

The Principles of Probability

From Formal Logic to Measure Theory
to the Principle of Indifference

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Contents

Preface	xi
1 Introduction	1
1.1 Deductive vs. inductive reasoning	1
1.2 The two sides of logic	2
1.3 The nature of probability	3
1.4 Potential areas of application	5
1.5 The principle of indifference	6
1.6 A philosophical aside	7
1.7 Constructing inductive logic	8
1.8 Inductive logic is natural	9
1.9 Outline of the book	10
2 Background	13
2.1 Ordinal and cardinal numbers	13
2.2 Boolean algebras	16
2.3 Measure spaces	17
2.4 Structures	21
2.5 Strings	24
3 Propositional Calculus	25
3.1 Formulas and deductive inference	27
3.2 Inductive statements and entire sets	35
3.3 Closed sets and inductive derivability	43
3.4 Connectivity, restrictions, and lifts	49
3.5 Generating inductive theories	58
4 Propositional Models	65
4.1 Models and deductive semantics	67
4.2 Inductive semantics	76
4.3 Counterexamples and resolutions I	86
4.4 Counterexamples and resolutions II	91
4.5 Independence	95

5	Predicate Logic	105
5.1	The syntax of predicate formulas	106
5.2	Predicate calculus	113
5.3	Predicate models	124
5.4	Predicate models and random variables	134
6	Real inductive theories	143
6.1	Definitorial extensions	145
6.2	Zermelo–Fraenkel set theory	152
6.3	Real inductive theories in ZFC ₋	158
6.4	Real inductive theories in ZFC	164
6.5	Limit theorems	171
6.6	Probabilities of probabilities	175
7	Principle of Indifference	181
7.1	Formulating the principle	181
7.2	Examples with a single object	187
7.3	Examples with multiple objects	192
7.4	Indifference and exchangeability	199
7.5	Examples on an interval	201
7.6	Examples in the plane	210
	Bibliography	221
	Index of Terms	223
	Index of Symbols	227

Detailed Contents

Preface	xi
1 Introduction	1
1.1 Deductive vs. inductive reasoning	1
1.2 The two sides of logic	2
1.3 The nature of probability	3
1.4 Potential areas of application	5
1.5 The principle of indifference	6
1.6 A philosophical aside	7
1.7 Constructing inductive logic	8
1.8 Inductive logic is natural	9
1.9 Outline of the book	10
2 Background	13
2.1 Ordinal and cardinal numbers	13
2.1.1 Ordinal numbers	13
2.1.2 Transfinite induction and recursion	14
2.1.3 Ordinal arithmetic	15
2.1.4 Cardinal numbers	15
2.2 Boolean algebras	16
2.3 Measure spaces	17
2.3.1 Generating σ -algebras	18
2.3.2 Complete measure spaces	18
2.3.3 Dynkin systems	19
2.3.4 Measurable functions and pushforwards	19
2.3.5 Measure space isomorphisms	20
2.4 Structures	21
2.4.1 Structure homomorphisms	22
2.4.2 The standard structure of arithmetic	22
2.4.3 Factor structures	23
2.4.4 Direct products of structures	23
2.5 Strings	24

3	Propositional Calculus	25
3.1	Formulas and deductive inference	27
3.1.1	Propositional formulas	27
3.1.2	A calculus of natural deduction	29
3.1.3	A Hilbert-type calculus	31
3.1.4	Deductive theories and logical equivalence	33
3.2	Inductive statements and entire sets	35
3.2.1	Seven of nine	36
3.2.2	Relative negation and certainty	37
3.2.3	Inductive vs. deductive inference	39
3.2.4	Generalizations of the addition rule	40
3.2.5	Generalizations of the multiplication rule	41
3.2.6	Generalizations of the continuity rule	42
3.3	Closed sets and inductive derivability	43
3.3.1	The rule of inductive extension	43
3.3.2	The rule of deductive extension	45
3.3.3	Pre-theories	45
3.3.4	Inductive theories	47
3.3.5	Inductive derivability	48
3.4	Connectivity, restrictions, and lifts	49
3.4.1	Connectivity properties	49
3.4.2	Connectivity and inductive inference	50
3.4.3	Restrictions	51
3.4.4	The lift of a pre-theory	52
3.4.5	Identifying lifts with inductive theories	54
3.4.6	Characterizing inductive theories	57
3.5	Generating inductive theories	58
3.5.1	Strongly connected equivalence	58
3.5.2	Intersections of inductive sets	60
3.5.3	A converse to the rule of logical implication	61
3.5.4	Inductive conditions	62
4	Propositional Models	65
4.1	Models and deductive semantics	67
4.1.1	Truth assignments	67
4.1.2	Strict models and Boolean functions	69
4.1.3	Models and satisfiability	70
4.1.4	Deductive consequence and soundness	71
4.1.5	Karp's completeness theorem	73
4.1.6	Inductive theories and Dynkin systems	74
4.1.7	The full completeness theorem	74
4.2	Inductive semantics	76
4.2.1	Inductive satisfiability	76
4.2.2	Models determine theories	76
4.2.3	Theories determine models	79
4.2.4	Consistency and satisfiability	80

4.2.5	Inductive consequence and completeness	81
4.2.6	Differing roots	83
4.2.7	The semantics of inductive conditions	84
4.3	Counterexamples and resolutions I	86
4.3.1	Every probability space is a model	86
4.3.2	Dynkin spaces	87
4.3.3	Entirety is not enough	89
4.4	Counterexamples and resolutions II	91
4.4.1	An unknown false statement	92
4.4.2	Karp's counterexamples	93
4.5	Independence	95
4.5.1	Dialog sets	96
4.5.2	Independence of two formulas	97
4.5.3	Independence of a sequence of formulas	98
4.5.4	A semantic characterization of independence	99
4.5.5	Fair coin flips	100
4.5.6	Biased coin flips	102
5	Predicate Logic	105
5.1	The syntax of predicate formulas	106
5.1.1	The alphabet and terms	107
5.1.2	Formulas	108
5.1.3	Formula induction and recursion	110
5.1.4	Variables and symbols	110
5.1.5	Substitutions	112
5.2	Predicate calculus	113
5.2.1	Free substitutions	114
5.2.2	Natural deduction	115
5.2.3	Constant expansions	117
5.2.4	Deduction with sentences	119
5.2.5	Tautologies and consistency	121
5.2.6	Deductive and inductive theories	121
5.2.7	Karp's calculus	122
5.3	Predicate models	124
5.3.1	Strict satisfiability	124
5.3.2	Models and deductive satisfiability	126
5.3.3	Deductive consequence and soundness	128
5.3.4	Deductive completeness	130
5.3.5	Peano arithmetic	132
5.3.6	Inductive consequence and completeness	134
5.4	Predicate models and random variables	134
5.4.1	Random variables as extralogical symbols	135
5.4.2	Extralogical symbols as functions	137
5.4.3	The relativity of randomness	138
5.4.4	Frames of reference	139
5.4.5	The natural frame of reference	141

6	Real inductive theories	143
6.1	Definitorial extensions	145
6.1.1	Defining individual symbols	145
6.1.2	Defining multiple symbols	146
6.1.3	Extensions and models	146
6.1.4	Deductive elimination	148
6.1.5	Inductive elimination	149
6.1.6	Primitive vs. defined symbols	151
6.2	Zermelo–Fraenkel set theory	152
6.2.1	Extensionality, union, and power set	152
6.2.2	Axiom schema of separation	152
6.2.3	Axiom schema of replacement	153
6.2.4	Definitorial extensions and shorthand	154
6.2.5	Axioms of infinity, foundation, and choice	154
6.2.6	Finitary and infinitary ZFC	155
6.2.7	Consistency of ZFC	155
6.3	Real inductive theories in ZFC_*	158
6.3.1	The set of natural numbers	158
6.3.2	Arithmetic operations	159
6.3.3	Peano arithmetic and nonstandard numbers	160
6.3.4	Real numbers in ZFC_*	161
6.3.5	The standard real structure	162
6.3.6	Embedding random variables in ZFC_*	163
6.4	Real inductive theories in ZFC	164
6.4.1	Real numbers and Borel sets	164
6.4.2	The real frame of reference	165
6.4.3	Sequences and limits	168
6.4.4	Measurable functions	169
6.4.5	Embedding random variables in ZFC	169
6.5	Limit theorems	171
6.5.1	Borel terms	171
6.5.2	Jointly Borel terms	171
6.5.3	Independence of terms	172
6.5.4	The law of large numbers for terms	173
6.5.5	The central limit theorem for terms	174
6.6	Probabilities of probabilities	175
6.6.1	Conditioning on terms	175
6.6.2	Versions of distributions	177
6.6.3	Indicator terms	178
6.6.4	Conditional expectation	179
6.6.5	The law of total probability	180

7 Principle of Indifference	181
7.1 Formulating the principle	181
7.1.1 Signature permutations	182
7.1.2 Deductive indifference	183
7.1.3 Inductive indifference	185
7.1.4 Structures, models, and indifference	186
7.2 Examples with a single object	187
7.2.1 Either it's true or it isn't	187
7.2.2 A single coin flip	188
7.2.3 A single trial	189
7.2.4 Success is good	190
7.2.5 Goodness is independent	191
7.2.6 Lowering the root	192
7.3 Examples with multiple objects	192
7.3.1 Three balls, two colors	193
7.3.2 Two balls, two colors	194
7.3.3 Random numbers	195
7.3.4 Random numbers and definitions	196
7.4 Indifference and exchangeability	199
7.4.1 Permutations of real inductive theories	199
7.4.2 Exchangeability	200
7.5 Examples on an interval	201
7.5.1 The interval $[0, 1]$	201
7.5.2 A point on a rod	202
Introduction	202
Notation in ZFC	203
Extralogical symbols and assumptions	203
Inductive hypotheses and conclusion	204
A model for the rod	204
Narrowing down the permutations	205
Proof of main result	208
7.5.3 Adding a defined constant	208
7.6 Examples in the plane	210
7.6.1 A point on a circle	211
7.6.2 Bertrand's paradox	214
Introduction	214
Notation in ZFC	215
A first pass at setting up the problem	215
The complete setup and conclusion	216
Bibliography	221
Index of Terms	223
Index of Symbols	227

Preface

In classical logic, we formalize deductive arguments. The conclusion of a deductive argument is known with certainty, provided its premises are true. Inductive arguments, on the other hand, are those whose conclusions are known only with some degree of plausibility. An argument in a courtroom, for example, is inductive. Its conclusion, at best, is only known “beyond a reasonable doubt.”

We present herein a formal system of inductive logic. The system contains deductive logic as a special case. It also uses a language that allows for countable conjunctions. When we restrict our attention to deductive arguments using only finite conjunctions, we obtain ordinary first-order logic. That is, first-order logic is embedded within this system. In particular, the system is capable of expressing the usual set theory of Zermelo and Fraenkel, and as such, can express any statement of modern mathematics.

This system of inductive logic gives rise to a probability calculus that is in complete agreement with modern, mathematical probability theory. In particular, the inductive statements in the formal language of this system can be interpreted in probability spaces. Moreover, any probability space, together with any collection of random variables, can be mapped in a natural way to such an interpretive model.

Inductive logic, however, is more expressive than ordinary probability theory. There are probabilistic ideas that are expressible in this system which cannot be formulated using only probability spaces and random variables. An example of such an idea is the principle of indifference, a heuristic notion originating with Laplace. Roughly speaking, it says that if we are “equally ignorant” about two possibilities, then we should assign them the same probability. The principle of indifference has no rigorous formulation in ordinary probability theory. It exists only as a heuristic. Moreover, its use has a history of being problematic and prone to apparent paradoxes. In our system of inductive logic, however, we provide a rigorous formulation of this principle, and illustrate its use through a number of typical examples.

The material herein makes use of (mostly) basic facts from mathematical logic and measure theory. We assume the reader is already familiar with the fundamentals of measure-theoretic probability theory. On the other hand, an effort has been made to accommodate readers with no familiarity in logic. The logical notions that we use are presented in a way that is mostly self-contained. Where this is not possible, explicit references to the literature are provided.

Chapter 1

Introduction

Strictly speaking, all our knowledge outside mathematics and demonstrative logic (which is, in fact, a branch of mathematics) consists of conjectures.

—George Polya, 1954 [26]

1.1 Deductive vs. inductive reasoning

Newton’s laws of motion were conjectures that Einstein showed us were false. But Einstein’s theory of relativity is also a conjecture that may one day be overturned. Physical laws, in general, are all conjecture. Each experimental confirmation makes them more plausible, but they can never be established with complete certainty.

In the courtroom, the prosecuting attorney must prove the guilt of the defendant, not with complete certainty, but only beyond a reasonable doubt. If the law required the prosecutor to achieve complete certainty, then everyone would be acquitted, because this would be impossible. Strictly speaking, then, we imprison people on the basis of conjecture.

Likewise, the conclusions of the historian, the economist, the chemist, and the medical researcher are all conjecture. None of them can establish their results with complete certainty.

And yet, many of these results have been shown to be so plausible that no one seriously doubts them. To borrow an example from Laplace [22], the sun will rise tomorrow. I cannot say this with complete certainty, but the degree of plausibility of this fact is so high, that my human mind cannot even perceive the sliver of doubt that is there.

The process by which these conjectures obtain varying degrees of plausibility is not the logic of the mathematician. A mathematical proof either establishes complete certainty, or it fails to say anything at all. The logic of mathematics is *deductive reasoning*. The logic of everything else is *inductive reasoning*, or as Polya calls it in [26], “plausible reasoning.”

Deductive reasoning is governed by rules. Over the course of human history, we have uncovered and formalized these rules, and today we have complete systems of deductive logic that codify this form of argumentation. The study of deductive logic, or mathematical logic, is presently a mature and sophisticated subdiscipline of mathematics that has had profound impacts in areas ranging from computer science to philosophy to mathematics itself.

Inductive reasoning also has rules. For example, if Hypotheses A implies Hypothesis B , and Hypothesis B is found to be true, then the plausibility of Hypothesis A increases. This rule is the basis for empirical science. If the theory of relativity predicts something about Mercury, and the prediction is confirmed by experiment, then the theory of relativity is made more plausible by this discovery. In [27], Polya calls this the “fundamental inductive pattern.”

But unlike its deductive counterpart, inductive reasoning has not been formalized in any universally accepted way. Since the time of Laplace, probability theory has seemed like the most promising candidate to formalize inductive reasoning. For example, in probability theory it is a provable fact that if A implies B , and if $P(A)$ and $P(B)$ are neither 0 nor 1, then $P(A | B) > P(A)$. This is a formalization of Polya’s fundamental inductive pattern.

Attempts to formalize inductive reasoning with probability can be traced at least back to Boole in 1854 [4]. But the recognition that we need some kind of probabilistic logic goes all the way back to Leibniz in the 17th century. Since then, mathematicians, philosophers, and physicists have all contributed to this endeavor. See, for instance, [19, 33, 29, 5, 21, 25, 13]. For a survey of the history of these efforts, see [11].

In the meantime, over the last 90 years, modern probability—by which we mean measure-theoretic probability theory—has grown into a powerful and hugely successful discipline. It started with Kolmogorov in 1933 [20] and, today, enjoys phenomenal success in all its areas of applications. Advances have been made in physics, finance, engineering, meteorology, telecommunications, biology, astronomy, artificial intelligence, and more. Time has shown us that if we want to do quantified inductive reasoning, there is no better tool than modern probability.

Perhaps, then, we should turn the effort on its head. As presented above, we have been regarding inductive reasoning as the fundamental concept, and probability as a tool by which to formalize it. Instead, we might consider modern probability to be the fundamental concept, and seek out the inductive logic that it represents. To clarify what this might mean, let us look to its analogue in deductive logic, or more specifically, first-order logic.

1.2 The two sides of logic

First-order logic is the usual logic of sentences that involve quantifiers and predicates. We may approach first-order logic in one of two ways. We may consider it through a syntactic calculus, or we may consider it through semantics and meaning. When viewed as a calculus, the central idea is the “proof.” A

proof is a sequence of sentences whose construction follows a given set of rules, and which terminates in the sentence that is being proved. A sentence φ is derivable from X , written $X \vdash \varphi$, if there exists a proof of φ from X .

When viewed via semantics, the central idea is the “structure.” A structure is a set with distinguished constants, functions, and relations. A group is a structure with an identity and a group operation. A tree is a structure with a root and an edge relation. Even a committee of senators with a chair and two subcommittees is a structure. Formal sentences can be interpreted in a structure, and once interpreted, the sentence is either true or false. We say that φ is a consequence of X , written $X \models \varphi$ if φ is true in every structure where X is true.

A priori, these two approaches to first-order logic have nothing to do with one another. And yet, thanks to Gödel’s completeness theorem, we know that they exactly coincide. That is, $X \vdash \varphi$ if and only if $X \models \varphi$. In other words, first-order logic *is* the logic of structures. By analogy, if we want an inductive logic that *is* the logic of probability, then we should build a complete inductive logic that has modern, measure-theoretic probability as its semantics.

That is exactly what we will do in this manuscript. We will construct inductive logic so that, like first-order logic, it has a calculus and a semantics. The calculus will be based on nine rules of inductive inference. It will allow us to derive probabilities without any sample spaces or measure theory. All that matters in the calculus are the logical relations between the sentences, and the rules of inductive inference. On the semantic side, statements are interpreted in what we call an “inductive model,” which is a probability measure on a set of structures. We will establish completeness, showing that the calculus and the semantics coincide. And we will show that the whole of measure-theoretic probability theory is properly embedded in the semantic side of inductive logic. That is, any probability space, together with a set of its random variables, can be mapped to an inductive model in a way that gives each outcome, event, and random variable a logical interpretation.

1.3 The nature of probability

What, then, is probability? Or more precisely, what body of ideas should the word “probability” refer to? As a mathematician, it is tempting to say that probability is simply measure-theoretic probability, the branch of mathematics built on Kolmogorov’s formalism. But to paraphrase Terence Tao [31], probability spaces are used to model probabilistic concepts. They are not the concepts themselves. As such, and in light of the work done here, we might say that probability is the logic of inductive reasoning. It is an abstract mode of logical reasoning that is reflected in two parallel systems: a calculus of inductive inference, and a semantic system of interpretation. It is within this semantic system that measure-theoretic probability resides.

Another way to assess the nature of probability is to look at the people who study it and ask what they actually do. In other words, what is a probabilist?

The traditional view is simple. Probability is a branch of mathematics, and a probabilist is a mathematician who specializes in it. A probabilist is just one of many mathematical specialists, each of whom is classified according to the kinds of structures they study. A group theorist, for example, studies structures satisfying the axioms of group theory, that is, groups. A probabilist, therefore, studies structures satisfying Kolmogorov's axioms, that is, probability spaces.

This view, however, misses something important. A probabilist typically specializes in a particular class of random variables and stochastic processes. They do not simply study probability spaces. A person who studies pure probability spaces, without any random variables, would be better described as a measure theorist. A probabilist, on the other hand, studies what we might call "modern probability models," which are probability spaces equipped with a collection of random variables. As noted above, every modern probability model is an inductive model. Hence, probabilists study inductive models, and any given probabilist will study a particular class of inductive models.

To clarify the situation, let us note that probability spaces are to sets as inductive models are to structures. That is, we have four kinds of object (sets, structures, probability spaces, and inductive models) that all stand in relation to one another. The simplest object is the set. We may extend the notion of the set in two directions. On the one hand, if we add a probability measure to it, we obtain a probability space. On the other hand, if we add constants, functions, and relations, we obtain a structure, which we use to interpret deductive logic. An inductive model, which we use to interpret inductive logic, can be obtained from either a structure or a probability space. By definition, if we take a set of structures and add a probability measure, we have an inductive model. Or we can start with a probability space and add a particular collection of random variables. In doing so, we obtain a modern probability model, which is embedded in the collection of inductive models.

These four kinds of objects are studied by different kinds of mathematicians. Sets are studied by set theorists. Probability spaces, without any random variables, are studied by measure theorists. Structures come in many varieties. For example, graphs are studied by graph theorists. Likewise, inductive models come in many varieties. An example is random graphs. A person who studies random graphs would be called a probabilist.

But a probabilist who studies random graphs is as much a graph theorist as they are a probabilist. Just as mathematicians are categorized according to the kinds of structures they study, probabilists are categorized according to the kinds of inductive models they study. There are probabilists who study stochastic PDEs, random matrices, stochastic control theory, and so on. And like their deterministic counterparts, any given probabilist will typically only specialize in one such area. The word "mathematician" is an umbrella term that includes many different areas. Likewise, the word "probabilist" is also an umbrella term that includes all these same areas, but seen through the lens of probability and inductive reasoning. The picture that emerges from this view is that probability is not just a branch of mathematics. It is a different logical paradigm with which to study other branches of mathematics.

1.4 Potential areas of application

As an extension of formal deductive logic, inductive logic has many potential areas of application. Computer science, for instance, is rooted in mathematical logic. Hence, any probabilistic extension of computer science is, in some way, connected to a probabilistic extension of logic. It follows that inductive logic could be relevant to any such areas. Examples might include quantum computing and artificial intelligence.

Inductive logic could also be applicable to philosophy. It is connected to the philosophy of science through Polya's fundamental inductive pattern. It is also highly relevant to the philosophical interpretations of probability, and through these, to epistemology. For instance, we can use inductive logic to formalize the principle of indifference. This principle is the heuristic notion that we ought to assign equal probabilities to sentences about which we are equally ignorant. This principle is intuitively self-evident, but historically problematic. It leads to apparent paradoxes and is without any mathematically rigorous formulation. We will have more to say about the principle of indifference later.

Formal deductive logic can also be used to analyze philosophical arguments. As such, inductive logic can be used to analyze philosophical arguments that are probabilistic in nature. Examples include the doomsday argument, the simulation hypothesis, responses to the so-called Sleeping Beauty problem, and arguments surrounding superintelligence and the technological singularity.

Inductive logic can be applied to mathematics itself, as well as statistics. Its relevance to probability and Bayesian statistics is obvious. Outside of probability, it is relevant wherever probabilistic methods are used. For example, in graph theory and combinatorics, probabilistic methods are often used to establish existence theorems and asymptotic results. At the foundations of mathematics, it offers us a new tool for working with undecidable sentences. Given a set of axioms and an undecidable sentence which they can neither prove nor disprove, deductive logic allow us to explore the theories obtained by either including or excluding that sentence from our axioms. With inductive logic, we may choose to postulate a probability for the undecidable sentence, and explore the probabilistic statements that follow from this assumption.

There are many potential areas of applications in physics. One obvious candidate is quantum mechanics, particularly its interpretations. The interpretations of quantum mechanics, such as the many-worlds and the Bohmian interpretations, cannot be decided upon through experiment, since all of their predictions are based on the common mathematical framework of quantum mechanics. The problem of deciding which one is correct or most useful is a philosophical problem. Inductive logic gives us a framework for axiomatizing these interpretations, and thereby more thoroughly analyzing their structure and consequences.

Besides quantum mechanics, inductive logic might also be applicable to statistical mechanics. This is particularly true since statistical mechanics is rooted in the principle of indifference. Its fundamental postulate is that, a priori, the microstates of an isolated system occur with equal probability.

Inductive logic, therefore, through the principle of indifference, touches upon the foundations of statistical mechanics.

The study of noisy dynamical systems is an area of physics where modern probability has found great success. In applications, the use of stochastic ordinary and partial differential equations is widespread. On the other hand, we can axiomatize classical (deterministic) mechanics in an extension of first-order logic that allows countable conjunctions. This extended language is exactly the one in which we will build inductive logic. We can therefore add probability and uncertainty directly into such an axiomatization. Does doing so lead us to the modern theory of stochastic dynamical systems? If not, how is it related to that theory? Determining this could provide insight into both inductive logic and the modern usage of stochastic differential equations.

1.5 The principle of indifference

(S:intro-PoI) The principle of indifference is the heuristic idea that if we are equally ignorant about two propositions, then we ought to assign them the same probability. This idea originated with Laplace, and is at the heart of what is now called the classical interpretation of probability.

The theory of chance consists in reducing all the events of the same kind to a certain number of cases equally possible, that is to say, to such as we may be equally undecided about in regard to their existence, and in determining the number of cases favorable to the event whose probability is sought. The ratio of this number to that of all the cases possible is the measure of this probability.

—Pierre-Simon Laplace, 1814 [22]

One of the most famous descriptions of the principle is due to Keynes.

The Principle of Indifference asserts that if there is no *known* reason for predicating of our subject one rather than another of several alternatives, then relatively to such knowledge the assertions of each of these alternatives have an *equal* probability. Thus *equal* probabilities must be assigned to each of several arguments, if there is an absence of positive ground for assigning *unequal* ones.

This rule, as it stands, may lead to paradoxical and even contradictory conclusions.

—John Maynard Keynes, 1921 [19]

It is difficult to overstate the importance of the principle of indifference. Not only is it at the heart of major areas of science, such as statistical mechanics. It is also central to our understanding of what probability is. Philosophically, it forms the basis of the classical interpretation of probability. But beyond philosophy, it is the everyday intuition of the common person as to why a

balanced die is fair or why it is important to thoroughly shuffle a deck of cards. And yet, as Keynes rightly points out, it has a problematic history. Even rather elementary applications of this principle can quickly lead to nonsensical results and apparent paradoxes. This is due, in no small part, to the fact that, for centuries, the principle of indifference has eluded attempts to make it rigorous. Without rigor, there are no precise conditions that can tell us whether our attempts to use it are legitimate. A notable formulation of the principle is given by Edwin T. Jaynes's [13], in a book posthumously published in 2003. He used this formulation as the basis of his maximum entropy principle, which plays a key role in statistical mechanics. But even Jaynes's formulation is non-rigorous. Moreover, there is no formulation of this principle in modern, measure-theoretic probability theory.

Within inductive logic, however, we will be able to formulate the principle of indifference. We will show that our formulation is a faithful representation of the principle. It is not simply an ad hoc condition to which we affix the name. Being mathematically rigorous, our formulation is as free from paradoxes as any proven mathematical theorem.

Moreover, we will see that the principle of indifference cannot be formulated using only the axioms of Kolmogorov. Its formulation requires the structure of inductive logic, both its syntactic structure and the semantic structures embedded in its models. As such, it exemplifies the fact that inductive logic is strictly broader than any theory of probability that is based on measure theory alone.

1.6 A philosophical aside

This book is, without question, a work of mathematics. It consists primarily of definitions, theorems, and proofs, with occasional intuitive prose to tie it together. On the other hand, it is hard not to see it as a work of philosophy, having something to say about the interpretation of probability.

According to [12], interpretations of probability generally address the questions:

- (1) What kinds of things, metaphysically, are probabilities?
- (2) What makes probability statements true or false?

To use inductive logic to answer these questions, we must first look to the definition of an inductive statement. An inductive statement is a triple, (X, φ, p) , where X is a set of sentences (in a formal language) called the "antecedent," φ is a sentence that we call the "consequent," and p is a real number satisfying $0 \leq p \leq 1$ which we call the "probability." Intuitively, we can think of (X, φ, p) as asserting that X partially entails φ , and that p is the degree of this partial entailment. To answer the first question, then, a probability is a relationship between X and φ . We might interpret this relationship as being logical, evidential, or purely subjective. Any such interpretation has no bearing on the inductive logical properties of (X, φ, p) .

Inductive logic is mostly unconcerned with the second question. Probabilities express relative likelihoods, given a set of sentences. We take it for granted that the sense of these likelihoods is understood. Our primary concern is the logical relationships between inductive statements. That is, how can we reason from hypotheses, which are themselves inductive statements, to an inductive conclusion.

It seems, then, that inductive logic hardly qualifies as an interpretation of probability. It does, however, make assumptions that rule out certain interpretations. It assumes that probabilities are relationships between sentences and that all probabilities are conditional. Inductive logic, therefore, is incompatible with any physical interpretations of probability, such as those based on frequencies or propensities. Beyond that, though, it appears to leave room for a range of evidential interpretations.

1.7 Constructing inductive logic

`<S:construct-sketch>` In this section, we give a big-picture overview of how inductive logic is constructed. It is assumed that the reader is already familiar with the basics of measure-theoretic probability theory.

The set of sentences and formulas that we consider is one that allows for countable conjunctions and disjunctions. This set is usually denoted in the literature by $\mathcal{L}_{\omega_1, \omega}$, though we denote it simply by \mathcal{L} . In the language \mathcal{L} , there is a well-understood deductive calculus (see, for instance, [16, 18]). That is, there is a well-established derivability relation \vdash , where $X \vdash \varphi$ indicates that $\varphi \in \mathcal{L}$ can be derived from $X \subseteq \mathcal{L}$. Our first task is to extend \vdash to inductive statements. We do this by defining a set of rules for inductive inference, and writing $P \vdash (X, \varphi, p)$ to mean that we can use these rules of inference to derive (X, φ, p) from the set of inductive statements P . We also define an inductive theory, which is a set that is closed under inductive inference. Finally, we adopt familiar shorthand, writing $P(\varphi \mid X) = p$ to mean that $(X, \varphi, p) \in P$. We say that $P(\varphi \mid X)$ exists if $(X, \varphi, p) \in P$ for some p .

Let us refer to the left-hand side of the turnstile symbol, \vdash , as the premises of the derivation $P \vdash (X, \varphi, p)$. From what we have described so far, our premises can only include statements of the form $P(\varphi \mid X) = p$. We may wish, however, to include premises of the form $P(\varphi \mid X) > 0$ or $P(\varphi \wedge \psi \mid X) = P(\varphi \mid X)P(\psi \mid X)$. To this end, we generalize inductive derivability by defining inductive conditions, typically denoted by calligraphic letters such as \mathcal{C} . Inductive conditions formalize generic assumptions we might make about an inductive theory. After defining inductive conditions, we extend inductive derivability from $P \vdash (X, \varphi, p)$ to $\mathcal{C} \vdash (X, \varphi, p)$.

We then turn our attention to the semantics of inductive logic. More specifically, we define a relation, \models , called the consequence relation. The derivability relation, \vdash , is concerned only with the syntax of sentences, formulas, and inductive statements. The consequence relation, on the other hand, is concerned with their interpretations.

To give interpretations to sentences, we must introduce models. For us, a model, or an inductive model, is a probability space, $(\Omega, \Sigma, \mathbb{P})$, where Ω is a set of structures. A structure is a set, together with some distinguished constants, functions, and relations. For example, $(\mathbb{N}, 1, +, <)$ is a structure. Structures are used to interpret the sentences in \mathcal{L} . We write $\omega \models \varphi$ to mean that, in the structure ω , the interpretation of φ is a true statement. We can think of an inductive model as a weighted collection of structures, where the weights represent relative likelihoods.

Suppose $\mathcal{S} = (\Omega, \Sigma, \mathbb{P})$ is a model and φ is a sentence. Then we define the set $\varphi_\Omega = \{\omega \in \Omega \mid \omega \models \varphi\}$. We say that \mathcal{S} satisfies φ , written $\mathcal{S} \models \varphi$, if $\overline{\mathbb{P}}\varphi_\Omega = 1$, where $\overline{\mathbb{P}}$ is the measure-theoretic completion of \mathbb{P} . For sets of sentences, we take $\mathcal{S} \models X$ to mean that $\mathcal{S} \models \varphi$ for all $\varphi \in X$. We say that X and X' are semantically equivalent if, for every model \mathcal{S} , we have $\mathcal{S} \models X$ if and only if $\mathcal{S} \models X'$.

More generally, we write $\mathcal{S} \models (X, \varphi, p)$ to mean there exists a set of sentences Y and a sentence ψ such that $\mathcal{S} \models Y$, the sets $Y \cup \{\psi\}$ and X are semantically equivalent, and

$$\frac{\overline{\mathbb{P}}\varphi_\Omega \cap \psi_\Omega}{\overline{\mathbb{P}}\psi_\Omega} = p.$$

For sets of inductive statements, we take $\mathcal{S} \models P$ to mean $\mathcal{S} \models (X, \varphi, p)$ for all $(X, \varphi, p) \in P$. A set P is called satisfiable if $\mathcal{S} \models P$ for some model \mathcal{S} .

Having defined satisfiability, we turn to the consequence relation. We define $X \models \varphi$ to mean that $\mathcal{S} \models \varphi$ whenever $\mathcal{S} \models X$, and we define $P \models (X, \varphi, p)$ to mean that P and (X, φ, p) satisfy certain connectivity requirements and that $\mathcal{S} \models (X, \varphi, p)$ whenever $\mathcal{S} \models P$. We then proceed to prove soundness and completeness, meaning that the two relations, \vdash and \models , are identical. In other words, we prove that $X \vdash \varphi$ if and only if $X \models \varphi$, and also that $P \vdash (X, \varphi, p)$ if and only if $P \models (X, \varphi, p)$.

1.8 Inductive logic is natural

In mathematics, a definition is neither correct nor incorrect. It simply is. Nevertheless, we understand that there is a metamathematical sense in which a definition can be “right” or “wrong.” Does it capture the intuitive idea it alleges to capture? Does it “fit” well with existing notions? Does it provide a sense of mathematical “unity” or “elegance?”

In constructing inductive logic, we have made several choices. We chose to work in the infinitary language, $\mathcal{L}_{\omega_1, \omega}$. We chose the rules of inductive inference that characterize the derivability relation, \vdash . And we chose a particular semantic interpretation in order to arrive at the consequence relation, \models . Were these the “right” choices? Are they natural? Or are they ad hoc choices whose only purpose is to force the conclusions we are aiming for? These are metamathematical questions, and so, for the most part, their answers are left to the judgment of the reader. But we certainly believe they are natural, and present here two comments which are related to this question.

The first comment concerns the rules of inductive inference used in defining \vdash . There are rules that govern probabilities 0 and 1, and relate them to deductive inference. There is a continuity rule, needed only in infinite systems. Otherwise, the two main operative rules are the familiar ones from elementary probability: the addition rule and the multiplication rule. We make no effort in this work to justify their adoption. Inductive logic may be seen as simply investigating their consequences, or as implying that they are self-evident. That said, in the paragraph after next, we will make one comment in their defense. Beyond that, the interested reader can consult [8] for an attempt to justify these elementary principles.

The second comment concerns \models . We essentially have four interconnected relations: deductive \vdash , inductive \vdash , deductive \models , and inductive \models . For deductive \models , we prove both completeness ($X \models \varphi$ implies $X \vdash \varphi$) and σ -compactness ($X \models \varphi$ implies $X_0 \models \varphi$ for some countable $X_0 \subseteq X$). As mentioned earlier, deductive \vdash is the usual relation used in $\mathcal{L}_{\omega_1, \omega}$. Deductive \models , however, is not. The usual consequence relation in $\mathcal{L}_{\omega_1, \omega}$ is the more straightforward $X \equiv \varphi$, meaning $\omega \equiv X$ implies $\omega \equiv \varphi$. It is well-known that \equiv is neither complete, nor σ -compact. For this reason, one might regard $\mathcal{L}_{\omega_1, \omega}$ as being deficient in some important ways. Our primary purpose for introducing \models is to model inductive reasoning. As an unintended but welcome consequence, we find that \models corrects these deficiencies. In this sense, then, we might view \models as the “right” semantic notion for $\mathcal{L}_{\omega_1, \omega}$, and also see $\mathcal{L}_{\omega_1, \omega}$ as the “right” language in which to build a probabilistic logic.

Turning this argument on its head, if deductive \models is the “right” semantic relation to pair with deductive \vdash , and inductive \models follows from deductive \models , then our definition of inductive \vdash is forced on us, if we want $P \models (X, \varphi, p)$ and $P \vdash (X, \varphi, p)$ to be equivalent. In other words, we are compelled to adopt the addition and multiplication rules, as well as all the other rules in the definition of \vdash .

1.9 Outline of the book

(S:outline) After presenting background material in Chapter 2, we proceed to the construction of inductive logic. In Chapters 3 and 4, we follow the sketch presented in Section 1.7 to construct a trimmed-down version of inductive logic in a propositional language without quantifiers or variables. The propositional version is capable of representing any probability space, but it does not explicitly represent any random variables.

In Chapter 5, we repeat the construction for the predicate language \mathcal{L} . Here, we are able to establish the connection between inductive logic and random variables. In the formal language \mathcal{L} , constants, functions, and relations are represented by what are called extralogical symbols. If s is an extralogical symbol and ω is a structure, then we can identify s with an actual constant, function, or relation in ω , which we denote by s^ω . The map $\omega \mapsto s^\omega$ is the starting point for connecting inductive logic to random variables.

Almost immediately, though, we face an obstacle. In probability theory, we are accustomed to having two very distinct kinds of objects: random variables and constants. In $\mathbb{P}\{X > 0\}$, for example, we can be quite certain that only thing random is X . In the formal language, \mathcal{L} , however, X , $>$, and 0 are just symbols. When we interpret them in a structure, we are faced with $\mathbb{P}\{X^\omega >^\omega 0^\omega\}$. Hence, not only might 0 be random, the inequality relation itself could be random!

These considerations lead us to an idea we call the relativity of randomness. To illustrate this, consider a simple system representing a coin flip. We might have just three extralogical symbols, h , t , and c , where we think of h and t as the sides of the coin, and c as the result. To be sure, there are models that match our intuition. That is, there are models in which h^ω and t^ω are fixed, and c^ω is random. But there are also models in which h^ω and t^ω are random, while c^ω is fixed. There are even models in which all three are random. In general, using model isomorphisms, we can force certain sets of extralogical symbols to be nonrandom. There are limits though. In the coin-flip example, there are no models in which all three symbols are fixed.

We show that in any inductive theory rich enough to contain the natural numbers, we can use model isomorphisms to force the natural numbers to be nonrandom. We call this the natural frame of reference.

Chapter 6 is concerned with constructing inductive theories that make statements about real numbers. We do this by constructing the real numbers, in the usual way, in axiomatic set theory. We are then able to make formal inductive statements about not only real numbers, but about any mathematical objects whatsoever. In this context, we illustrate how all the usual results of probability theory can be formulated in inductive logic. This includes the law of large numbers, the central limit theorem, and conditional expectation.

The principle of indifference is the topic of Chapter 7. We formulate it as an inductive condition. To describe it, let L be the set of extralogical symbols in our language and consider a bijection $\pi : L \rightarrow L$ that preserves the type and arity of the symbols. For instance, if r is a binary relation symbol, then so is r^π . Given a sentence φ , let φ^π be the sentence obtained from φ by replacing every instance of s with s^π . Similarly, let $X^\pi = \{\varphi^\pi \mid \varphi \in X\}$. We say that X is invariant under π if X and X^π are logically equivalent.

In the deductive calculus, if we take a proof of φ from X and transform each step of that proof by π , then we obtain a proof of φ^π from X^π . In other words, $X \vdash \varphi$ if and only if $X^\pi \vdash \varphi^\pi$. The principle of indifference is the natural extension of this to the inductive calculus.

Let P be an inductive theory. Then P satisfies the principle of indifference if $P(\varphi^\pi \mid X^\pi) = P(\varphi \mid X)$. In particular, if X is invariant under π , then, given X , the sentences φ and φ^π should be assigned the same probability.

After formulating the principle of indifference, we present several examples, beginning with simple discrete examples and continuing to examples involving intervals and circles. Our final example is an analysis of Bertrand's paradox, a famous counterintuitive illustration of the principle of indifference.

Chapter 2

Background

(Ch:background) It is assumed that the reader is familiar with the basics of measure-theoretic probability theory. In Section 2.3, we introduce some basic concepts from measure theory. The reader should already be familiar with almost everything in that section. We introduce them only to establish notation, cite some lesser known results, and establish new definitions that we will need later.

Familiarity with mathematical logic would be helpful but is not required. The concepts from logic that we use are all presented as we need them, in a mostly self-contained way.

We will adopt the convention that 0 is a natural number. However, we will use the notation $\mathbb{N} = \{1, 2, 3, \dots\}$ and $\mathbb{N}_0 = \{0, 1, 2, 3, \dots\}$. In other words, \mathbb{N} is the set of positive integers, and \mathbb{N}_0 is the set of natural numbers.

2.1 Ordinal and cardinal numbers

2.1.1 Ordinal numbers

(S:ordinals) For further details about ordinal and cardinal numbers, see any basic text on set theory, such as [6] or [9].

An *ordinal (number)* is a well-ordered set, $(\alpha, <)$, such that if $x \in \alpha$, then $x = \{y \in \alpha : y < x\}$. If α is an ordinal and $x, y \in \alpha$, then $x < y$ if and only if $x \in y$ if and only if $x \subset y$. (Here, we use \subseteq for subset and \subset for proper subset.) Moreover, every $x \in \alpha$ is itself an ordinal. Thus, an ordinal is a set of ordinals that is well-ordered by \in .

The collection of all ordinals is not a set. (If it were, then it would be an ordinal that is an element of itself, which, as we note below, is impossible.) It is, nonetheless, “well-ordered” in a certain sense. More specifically, if α and β are ordinals, then $\alpha \in \beta$ if and only if $\alpha \subset \beta$, so that \in has the properties of a strict partial order. Also, for any two ordinals, α and β , exactly one of the following is true: $\alpha \in \beta$, $\alpha = \beta$, or $\beta \in \alpha$. If α and β are ordinals, we will write $\alpha < \beta$ for $\alpha \in \beta$. Finally, every nonempty collection of ordinals has a smallest element. This last fact can be made rigorous in axiomatic set theory

as a so-called “theorem schema,” but for our present purposes, we do not need such a formalization.

From these facts, it follows that there is a smallest ordinal. The smallest ordinal is \emptyset , since there is no ordinal with $\alpha \in \emptyset$. The collection of nonempty ordinals has a smallest element, meaning there is a second smallest ordinal. The second smallest ordinal is $\{\emptyset\}$, since there is no ordinal with $\emptyset \subset \alpha \subset \{\emptyset\}$.

If α is an ordinal, let $s(\alpha)$ denote the ordinal $\alpha \cup \{\alpha\}$. We call $s(\alpha)$ the successor of α . Note that there is no β with $\alpha \subset \beta \subset s(\alpha)$. Also note that $s(\emptyset) = \{\emptyset\}$. From here, we see that the third smallest ordinal is $s(\{\emptyset\}) = \{\emptyset, \{\emptyset\}\}$, the fourth smallest ordinal is $\{\emptyset, \{\emptyset\}, \{\emptyset, \{\emptyset\}\}$, and so on.

For each $n \in \mathbb{N}_0$, define the ordinal n' by setting $0' = \emptyset$ and $(n+1)' = s(n')$. Let $\mathbb{N}'_0 = \{n' \mid n \in \mathbb{N}_0\}$. The map $n \mapsto n'$ is a bijection from \mathbb{N}_0 to \mathbb{N}'_0 which preserves the natural ordering on \mathbb{N}_0 . That is, $m < n$ if and only if $m' < n'$.

It is straightforward to verify that \mathbb{N}'_0 itself is an ordinal. Since \mathbb{N}'_0 is an infinite set and every $\alpha \in \mathbb{N}'_0$ is a finite set, it follows that \mathbb{N}'_0 is the smallest infinite ordinal. The usual notation for \mathbb{N}'_0 is ω . We will use this notation sparingly, since it conflicts with the usual usage of ω in probability theory. The reader must generally rely on context to see the distinction, although to make things easier, we will use bold font when using ω to denote \mathbb{N}'_0 . As we will discuss further in Section 2.4.2, we will sometimes find it useful to identify n with n' , giving us a representation of each natural number as a set.

Every well-ordered set is order isomorphic to a unique ordinal. In particular, no two distinct ordinals are order isomorphic. If $\alpha = s(\beta)$ for some ordinal β , then α is called a *successor ordinal*. Otherwise, α is called a *limit ordinal*. An ordinal α is a limit ordinal if and only if $\alpha = \bigcup_{\xi < \alpha} \xi$, and it is a successor ordinal if and only if $\alpha = s(\bigcup_{\xi < \alpha} \xi)$. The smallest limit ordinal is $0'$ and the second smallest limit ordinal is ω .

If α is an ordinal, then an α -*sequence* is a function whose domain is α , typically denoted by $\langle x_\xi \mid \xi < \alpha \rangle$. Such a sequence is said to have length α . If $\alpha = n'$, then these are n -tuples. If $\alpha = \omega$, then these are ordinary sequences indexed by \mathbb{N}_0 .

2.1.2 Transfinite induction and recursion

Let α be an ordinal and let S be a set. The principle of transfinite induction says that if $0' \in S$, and $\beta \subset S$ implies $\beta \in S$ for all $\beta < \alpha$, then $\alpha \subseteq S$. The principle of transfinite recursion states that if $G : \bigcup_{\xi < \alpha} S^\xi \rightarrow S$, then there exists a unique sequence $\langle x_\xi \mid \xi < \alpha \rangle$ such that $x_\beta = G(\langle x_\xi \mid \xi < \beta \rangle)$ for all $\beta < \alpha$. In proofs that use transfinite induction or recursion, the inductive/recursive step is typically broken into cases according to whether α is a successor ordinal or a limit ordinal.

2.1.3 Ordinal arithmetic

We define ordinal addition recursively by

$$\begin{aligned}\alpha + 0' &= \alpha, \\ \alpha + s(\beta) &= s(\alpha + \beta), \text{ and} \\ \alpha + \beta &= \bigcup_{\xi < \beta} (\alpha + \xi), \text{ if } \beta > 0' \text{ is a limit ordinal.}\end{aligned}$$

Note that $\alpha + 1' = s(\alpha)$. Ordinal addition is associative, but not commutative. For instance, $1' + \omega = \omega$, but $\omega + 1' > \omega$. Ordinal subtraction on the left is always possible. That is, if $\alpha \leq \beta$, then there exists a unique ordinal γ such that $\alpha + \gamma = \beta$. Consequently, ordinal addition is left-cancellative, meaning that $\alpha + \beta = \alpha + \gamma$ implies $\beta = \gamma$. Moreover, it is strictly increasing in the right argument, meaning that $\alpha + \beta < \alpha + \gamma$ if and only if $\beta < \gamma$. These facts can be used, for instance, to prove that the function $f : \gamma \rightarrow (\beta + \gamma) \setminus \beta$ given by $f(\xi) = \beta + \xi$ is a bijection.

We define ordinal multiplication recursively by

$$\begin{aligned}\alpha \cdot 0' &= 0', \\ \alpha \cdot s(\beta) &= (\alpha \cdot \beta) + \alpha, \text{ and} \\ \alpha \cdot \beta &= \bigcup_{\xi < \beta} (\alpha \cdot \xi), \text{ if } \beta > 0' \text{ is a limit ordinal.}\end{aligned}$$

Ordinal multiplication is associative, but not commutative. For instance, $2' \cdot \omega = \omega$, but $\omega \cdot 2' = \omega + \omega > \omega$. Multiplication is left distributive over addition, but not right distributive. For instance, $1' \cdot \omega + 1' \cdot \omega = \omega + \omega > \omega$.

When restricted to \mathbb{N}'_0 , ordinal addition and multiplication agree with ordinary addition and multiplication. That is, $(m + n)' = m' + n'$ and $(m \cdot n)' = m' \cdot n'$.

2.1.4 Cardinal numbers

If X and Y are sets, we write $X \sim Y$ to mean there exists a bijection $f : X \rightarrow Y$. Every set can be well-ordered, and every well-ordered set is isomorphic to a unique ordinal. Thus, for every set X , there exists an ordinal α such that $X \sim \alpha$. This α , however, is not unique, since a set can be well-ordered in multiple ways.

We define the *cardinality* of a set X to be the smallest ordinal α such that $X \sim \alpha$, and we write $|X| = \alpha$ or $\text{card}(X) = \alpha$. A *cardinal (number)* is an ordinal α that is the cardinality of some set. Since $|\alpha + 1| = |\alpha|$ whenever α is infinite, it follows that every infinite cardinal number is a limit ordinal.

Every finite ordinal is a cardinal, and ω is a cardinal. No other countably infinite ordinal besides ω is a cardinal. Since uncountable sets exist, we know that uncountable cardinals exist. Let ω_1 be the first uncountable cardinal, which is also the first uncountable ordinal.

For any set X , we have $|\mathfrak{P}X| > |X|$, where $\mathfrak{P}X$ is the power set of X . Hence, for any cardinal number κ , we have $|\mathfrak{P}\kappa| > \kappa$, which means that the

collection of cardinal numbers greater than κ is nonempty. Let κ^+ be the smallest cardinal greater than κ . In particular, $\omega_1 = \omega^+$. Any cardinal that is of the form κ^+ for some κ is called a *successor cardinal*. Otherwise, it is called a *limit cardinal*.

Note that ω_1 is the set of all countable ordinals. Every countable sequence of countable ordinals has a countable upper bound. More specifically, if α is an ordinal with $\alpha < \omega_1$, and $\langle \beta_\xi \mid \xi < \alpha \rangle$ is an α -sequence of ordinals with $\beta_\xi < \omega_1$ for all ξ , then there exists an ordinal $\gamma < \omega_1$ such that $\beta_\xi \leq \gamma$ for all ξ . In fact, we may take $\gamma = \bigcup_{\xi < \alpha} \beta_\xi$. Since a countable union of countable sets is countable, it follows that $\gamma < \omega_1$.

More generally, a cardinal κ is called *regular* if, whenever $\alpha < \kappa$ and $\langle \beta_\xi \mid \xi < \alpha \rangle$ is an α -sequence of ordinals with $\beta_\xi < \kappa$ for all ξ , we have $\bigcup_{\xi < \alpha} \beta_\xi < \kappa$. Since a finite union of finite sets is finite, ω is regular. Since a countable union of countable sets is countable, ω_1 is regular.

A cardinal number κ is called a *strong limit cardinal* if, for all cardinals $\lambda < \kappa$, we have $|\mathfrak{P} \lambda| < \kappa$. Since $|\mathfrak{P} \kappa| \geq \kappa^+$, a strong limit cardinal is always a limit cardinal. Since the power set of a finite set is finite, ω is a strong limit cardinal. However, the power set of a countable set can be uncountable, so ω_1 is not a strong limit cardinal.

The cardinal ω is both regular and a strong limit cardinal. We cannot “reach” or “access” ω from the finite sets that lie below it, using the operations of union and power set. A cardinal κ is called *strongly inaccessible* if it is uncountable, regular, and a strong limit cardinal.

2.2 Boolean algebras

A *Boolean algebra* is a partially ordered set (B, \leq) such that

- (i) $x \vee y := \sup\{x, y\}$ exists for all $x, y \in B$,
- (ii) $x \wedge y := \inf\{x, y\}$ exists for all $x, y \in B$,
- (iii) $x \wedge (y \vee z) = (x \wedge y) \vee (x \wedge z)$ for all $x, y, z \in B$,
- (iv) $x \vee (y \wedge z) = (x \vee y) \wedge (x \vee z)$ for all $x, y, z \in B$,
- (v) there exists $0, 1 \in B$ such that $0 \leq x \leq 1$ for all $x \in B$, and
- (vi) for all $x \in B$, there exists $\neg x \in B$ such that $x \wedge \neg x = 0$ and $x \vee \neg x = 1$.

It can be shown that in a Boolean algebra, $\neg x$ is unique. If

- (vii) $\bigvee_{x \in C} x := \sup C$ exists for all countable $C \subseteq B$, and
- (viii) $\bigwedge_{x \in C} x := \inf C$ exists for all countable $C \subseteq B$,

then B is a *Boolean σ -algebra*. The smallest Boolean algebra is the *degenerate* Boolean algebra with only one element, in which $0 = 1$. The smallest

nondegenerate Boolean algebra is $\mathbf{B} = \{0, 1\}$ with the usual meaning of \leq . The Boolean algebra \mathbf{B} is also clearly a Boolean σ -algebra.

Let B be a Boolean σ -algebra and $N \subseteq B$. Then N is a σ -ideal of B if $0 \in N$, $\bigvee_{x \in C} x \in N$ whenever $C \subseteq N$ is countable, and $x \in N$ implies $y \in N$ for all $y \leq x$.

If N is a σ -ideal of B and $x, y \in B$, we say $x \equiv y \pmod{N}$ if $(x \wedge \neg y) \vee (\neg x \wedge y) \in N$. The Boolean operations of B determine Boolean operations on the set of equivalence classes, B/N , making B/N into a Boolean σ -algebra.

A *Boolean measure* on a Boolean σ -algebra B is a function $m : B \rightarrow [0, \infty]$ such that $m(x) = 0$ if and only if $x = 0$, and $m(\bigvee_n x_n) = \sum_n m(x_n)$ whenever $i \neq j$ implies $x_i \wedge x_j = 0$. A *Boolean measure space* is a pair (B, m) where B is a Boolean σ -algebra and m is a Boolean measure on B .

Let (B, m) and (B', m') be Boolean measure spaces. A *homomorphism* from (B, m) to (B', m') is a function $g : B \rightarrow B'$ such that $x \leq y$ implies $g(x) \leq g(y)$, and $m'(g(x)) = m(x)$. The function g is an *isomorphism* if it is a bijection. In that case, g^{-1} is an isomorphism from (B', m') to (B, m) . If an isomorphism from (B, m) to (B', m') exists, then we say that (B, m) and (B', m') are *isomorphic*.

2.3 Measure spaces

(S:meas-spaces) Let Ω be a nonempty set and Σ a collection of subsets of Ω . Then Σ is a σ -algebra (of sets) on Ω if Σ is nonempty and closed under complements and countable unions. In this case, we call (Ω, Σ) a *measurable space*. A set $A \in \Sigma$ is called a *measurable set*. Note that a σ -algebra is a Boolean σ -algebra when it is equipped with the partial order \subseteq .

The intersection of any family of σ -algebras is a σ -algebra. If \mathcal{E} is any collection of subsets of Ω , then $\sigma(\mathcal{E})$ denotes the smallest σ -algebra containing \mathcal{E} , and is called the σ -algebra generated by \mathcal{E} . A *measure* on (Ω, Σ) is a function $\mu : \Sigma \rightarrow [0, \infty]$ such that $\mu \emptyset = 0$ and $\mu \bigcup_1^\infty A_n = \sum_1^\infty \mu A_n$ whenever $\{A_n\} \subseteq \Sigma$ is a pairwise disjoint sequence of measurable sets. In this case, (Ω, Σ, μ) is a *measure space*. A *measure subspace* of (Ω, Σ, μ) is a measure space (Ω, Σ', ν) , where $\Sigma' \subseteq \Sigma$ and $\nu = \mu|_{\Sigma'}$. If $\mu \Omega < \infty$, then μ is a *finite measure*. If $\mu \Omega = 1$, then μ is a *probability measure* and (Ω, Σ, μ) is a *probability space*. Any measure subspace of a probability one is also a *probability space*. In a probability space, it is customary to call the elements $\omega \in \Omega$ *outcomes*, and the measurable sets $A \in \Sigma$ *events*.

Let (Ω, Σ, μ) be a measure space. A μ -null set (or just a *null set*) is a set $A \in \Sigma$ with $\mu A = 0$. The collection of all null sets is denoted by \mathcal{N}_μ . If A and B are subsets of Ω , then we write $A = B$ μ -almost everywhere, if $A \Delta B$ is a subset of a null set, where

$$A \Delta B = (A \cap B^c) \cup (A^c \cap B)$$

is the *symmetric difference*. We will usually abbreviate this as $A = B$ μ -a.e., or if the measure is understood, as just $A = B$ a.e. If μ is a probability measure,

then we instead write $A = B$ μ -almost surely, abbreviated as $A = B$ μ -a.s. or $A = B$ a.s. We also write $A \subseteq B$ a.e. if $A \cap B^c$ is a subset of a null set. More generally, if f and g are functions with domain Ω , then we write $f = g$ a.e. if there exists $N \in \mathcal{N}_\mu$ such that $f(\omega) = g(\omega)$ for all $\omega \in N^c$.

2.3.1 Generating σ -algebras

(S:gen-sig-alg) Let Ω be a nonempty set and \mathcal{E} a collection of subsets of Ω . Then $\sigma(\mathcal{E})$ can be constructed from \mathcal{E} in an iterative fashion using transfinite recursion. Let $\mathcal{E}_0 = \mathcal{E}$. For an ordinal $\alpha < \omega_1$, let

$$\mathcal{E}'_\alpha = \mathcal{E}_\alpha \cup \{V^c \mid V \in \mathcal{E}_\alpha\},$$

and

$$\mathcal{E}_{\alpha+1} = \mathcal{E}'_\alpha \cup \{\bigcap \mathcal{D} \mid \mathcal{D} \subseteq \mathcal{E}'_\alpha \text{ is nonempty and countable}\}.$$

Here, countable means finite or countably infinite. If α is a limit ordinal, define $\mathcal{E}_\alpha = \bigcup_{\xi < \alpha} \mathcal{E}_\xi$.

By transfinite induction, $\mathcal{E}_\alpha \subseteq \sigma(\mathcal{E})$ for all $\alpha < \omega_1$, so that $\bigcup_{\alpha < \omega_1} \mathcal{E}_\alpha \subseteq \sigma(\mathcal{E})$. Clearly, $\bigcup_{\alpha < \omega_1} \mathcal{E}_\alpha$ is nonempty and closed under complements. Since every countable sequence of countable ordinals has a countable upper bound, it is also closed under countable intersections. Therefore, $\bigcup_{\alpha < \omega_1} \mathcal{E}_\alpha$ is a σ -algebra containing \mathcal{E} , which gives $\sigma(\mathcal{E}) \subseteq \bigcup_{\alpha < \omega_1} \mathcal{E}_\alpha$. This shows that $\sigma(\mathcal{E}) = \bigcup_{\alpha < \omega_1} \mathcal{E}_\alpha$.

For each $V \in \sigma(\mathcal{E})$, we define the *rank of V (with respect to \mathcal{E})*, which we denote by $\text{rk } V$, to be the smallest $\alpha < \omega_1$ such that $V \in \mathcal{E}_\alpha$. Note that $\text{rk } V$ is always a successor ordinal.

2.3.2 Complete measure spaces

Given a measure μ on (Ω, Σ) , we define the associated *outer measure* by

$$\mu^*A = \inf\{\mu B : A \subseteq B \text{ and } B \in \Sigma\}.$$

Note that μ^*A is defined for every $A \subseteq \Omega$. Similarly, we define the *inner measure* by

$$\mu_*A = \sup\{\mu B : B \subseteq A \text{ and } B \in \Sigma\}.$$

For any $A \subseteq \Omega$, we have $\mu_*A \leq \mu^*A$.

A *negligible set* is a (not necessarily measurable) subset of a null set. A negligible set that is also measurable is necessarily a null set. A measure space is called *complete* if every negligible set is measurable. A probability space is complete if and only if every superset of a set of measure 1 is measurable. In a complete measure space, if A is measurable and $A = B$ a.e., then B is measurable.

If (Ω, Σ, μ) is a measure space, then

$$\overline{\Sigma} = \{A \cup B : A \in \Sigma \text{ and } B \text{ is negligible}\}$$

is a σ -algebra called the *completion of Σ with respect to μ* . There is a unique measure $\bar{\mu}$ on $(\Omega, \bar{\Sigma})$ that agrees with μ on Σ and makes $(\Omega, \bar{\Sigma}, \bar{\mu})$ into a complete measure space. The measure $\bar{\mu}$ is called the *completion of μ* . Note that a set is $\bar{\mu}$ -null if and only if it is a subset of a μ -null set. Also, if $A \subseteq \Omega$ and $\mu^*A < \infty$, then $A \in \bar{\Sigma}$ if and only if $\mu_*A = \mu^*A$ (see, for instance, [7, Proposition 1.5.5]).

Let (Ω, Σ, μ) be a measure space and let $A \subseteq \Omega$. Then

$$\sigma(\Sigma \cup \{A\}) = \{(B \cap A) \cup (C \cap A^c) : B, C \in \Sigma\}.$$

Suppose μ is a finite measure and let $\alpha \geq 0$ satisfy $\mu_*A \leq \alpha \leq \mu^*A$. Then there exists a measure ν on $(\Omega, \sigma(\Sigma \cup \{A\}))$ such that $\nu|_{\Sigma} = \mu$ and $\nu A = \alpha$ (see, for instance, [7, Exercise 1.5.12] or [3, Theorem 1.12.14]).

2.3.3 Dynkin systems

Let Ω be a nonempty set and Δ a collection of subsets of Ω . Then Δ is a *Dynkin system*, or *λ -system*, if

- (i) $\Omega \in \Delta$,
- (ii) if $A, B \in \Delta$ with $A \subseteq B$, then $B \setminus A \in \Delta$, and
- (iii) If $\{A_n\} \subseteq \Delta$ with $A_n \subseteq A_{n+1}$, then $\bigcup_n A_n \in \Delta$.

Equivalently, one can define Δ to be a Dynkin system if it is nonempty and satisfies

- (i)' if $A \in \Delta$, then $A^c \in \Delta$
- (ii)' if $\{A_n\} \subseteq \Delta$ are pairwise disjoint, then $\bigcup_n A_n \in \Delta$.

Every σ -algebra is a Dynkin system. Conversely, a Dynkin system that is closed under (finite) intersections is a σ -algebra.

The intersection of any family of Dynkin systems is a Dynkin system. If \mathcal{E} is any collection of subsets of Ω , then there is a smallest Dynkin system containing \mathcal{E} , called the *Dynkin system generated by \mathcal{E}* .

If \mathcal{E} is a collection of subsets of Ω , then \mathcal{E} is a *π -system* if $A, B \in \mathcal{E}$ implies $A \cap B \in \mathcal{E}$. Dynkin's π - λ theorem states that if Δ is a Dynkin system, \mathcal{E} is a π -system, and $\mathcal{E} \subseteq \Delta$, then $\sigma(\mathcal{E}) \subseteq \Delta$.

Let Δ be a Dynkin system on Ω and let $B \in \Delta$. Define $\Delta|_B = \{A \subseteq \Omega : A \cap B \in \Delta\}$. Then $\Delta|_B$ is a Dynkin system on Ω called the *restriction of Δ to B* . If Δ is a σ -algebra, then $\Delta|_B$ is a σ -algebra.

2.3.4 Measurable functions and pushforwards

Let (Ω, Σ, μ) be a measure space and (S, Γ) a measurable space. A function $f : \Omega \rightarrow S$ is *measurable*, or *(Σ, Γ) -measurable*, if $U \in \Gamma$ implies $f^{-1}(U) \in \Sigma$.

If $A \subseteq \Omega$, then 1_A denotes the *indicator function* of A , and is defined by

$$1_A(\omega) = \begin{cases} 1 & \text{if } \omega \in A, \\ 0 & \text{if } \omega \notin A. \end{cases}$$

The function 1_A can be regarded as taking values in the measurable space (S, Γ) , where $S = \{0, 1\}$ and $\Gamma = \mathfrak{P}S$. In that case, if $A \subseteq \Omega$, then A is measurable if and only if 1_A is measurable.

Suppose f is a measurable function from a measure space (Ω, Σ, μ) to a measurable space (S, Γ) . Since $f^{-1}(\emptyset) = \emptyset$, $f^{-1}(U^c) = f^{-1}(U)^c$, and $f^{-1}(\bigcup_n U_n) = \bigcup_n f^{-1}(U_n)$, it follows that $\mu \circ f^{-1}$ is a measure on (S, Γ) , called the *pushforward* of μ .

2.3.5 Measure space isomorphisms

Let (Ω, Σ, μ) be a measure space and \mathcal{N}_μ the collection of null sets. The set \mathcal{N}_μ is a σ -ideal of Σ , and $A = B$ a.e. if and only if $A = B \pmod{\mathcal{N}_\mu}$. The equivalence classes modulo \mathcal{N}_μ are $[A]_\mu = \{B \in \Sigma : A = B \text{ a.e.}\}$, and the set of equivalence classes, Σ/\mathcal{N}_μ , is a Boolean σ -algebra. If we define $m_\mu : \Sigma/\mathcal{N}_\mu \rightarrow [0, \infty]$ by $m_\mu([A]_\mu) = \mu(A)$, then m_μ is well-defined and is a Boolean measure on Σ/\mathcal{N}_μ . We call $(\Sigma/\mathcal{N}_\mu, m_\mu)$ the *Boolean measure space corresponding to* (Ω, Σ, μ) . We say that two measure spaces are *isomorphic* if their corresponding Boolean measure spaces are isomorphic.

Let (Ω, Σ, μ) and (S, Γ, ν) be measure spaces, let $f : \Omega \rightarrow S$ be measurable, and assume $\nu = \mu \circ f^{-1}$. Note that $\mu f^{-1}(U) \triangle f^{-1}(V) = \nu U \triangle V$. Hence, $U = V$ ν -a.e. if and only if $f^{-1}(U) = f^{-1}(V)$ μ -a.e. It follows, therefore, that f determines an injective function $g : \Gamma/\mathcal{N}_\nu \rightarrow \Sigma/\mathcal{N}_\mu$ given by $g([U]_\nu) = [f^{-1}(U)]_\mu$. This function is, in fact, a homomorphism from $(\Gamma/\mathcal{N}_\nu, m_\nu)$ to $(\Sigma/\mathcal{N}_\mu, m_\mu)$. Hence, g is an isomorphism if and only if it is surjective, that is, if and only if

$$[A]_\mu \in \Sigma/\mathcal{N}_\mu \text{ implies } g([U]_\nu) = [A]_\mu \text{ for some } [U]_\nu \in \Gamma/\mathcal{N}_\nu.$$

We can rewrite this in terms of f to state the following. For two measure spaces, (Ω, Σ, μ) and (S, Γ, ν) , to be isomorphic, it suffices that there exists a measurable function $f : \Omega \rightarrow S$ such that $\nu = \mu \circ f^{-1}$, and

$$\text{for all } A \in \Sigma, \text{ there exists } U \in \Gamma \text{ such that } f^{-1}(U) = A \text{ } \mu\text{-a.e.}$$

If such an f exists, we say that f *induces an isomorphism* from (Ω, Σ, μ) to (S, Γ, ν) . Note that in this case, f also induces an isomorphism from $(\Omega, \overline{\Sigma}, \overline{\mu})$ to $(S, \overline{\Gamma}, \overline{\nu})$.

Two measure spaces, (Ω, Σ, μ) and (S, Γ, ν) , are *pointwise isomorphic* if there exists a measurable bijection $f : \Omega \rightarrow S$ such that f^{-1} is measurable and $\nu = \mu \circ f^{-1}$. In that case, f is called a *pointwise isomorphism*. It is straightforward to verify that if f is a pointwise isomorphism from (Ω, Σ, μ) to (S, Γ, ν) , then it is also a pointwise isomorphism from $(\Omega, \overline{\Sigma}, \overline{\mu})$ to $(S, \overline{\Gamma}, \overline{\nu})$. In

other words, if two measure spaces are pointwise isomorphic, then so are their completions.

Note that pointwise isomorphic measure spaces are isomorphic, and a pointwise isomorphism induces an isomorphism. The *pointwise isomorphism class* of (Ω, Σ, μ) is the collection of all measure spaces that are pointwise isomorphic to (Ω, Σ, μ) .

Let (Ω, Σ, μ) be a measure space and let S be a set. Let $h : \Omega \rightarrow S$ be a function. Define $\Gamma = \{A \subseteq S \mid h^{-1}(A) \in \Sigma\}$ and define $\nu = \mu \circ h^{-1}$. Then (S, Γ, ν) is a measure space and h is a measurable function from Ω to S . We call (S, Γ, ν) the *measure space image* of (Ω, Σ, μ) under h . If h is a bijection, then h is a pointwise isomorphism from (Ω, Σ, μ) to (S, Γ, ν) .

2.4 Structures

Let A be a set. If f is a function with domain A and $a \in A$, then the value of f at a will be denoted variously by $f(a)$, fa or a^f . Note that \emptyset is a function. In fact, \emptyset is the unique function with domain \emptyset .

If $n \in \mathbb{N}$, then A^n is the set of n -tuples, $\vec{a} = \langle a_1, \dots, a_n \rangle$. We let $A^0 = \{\emptyset\}$. If f is a function with domain A^n , we write $f\vec{a} = f(a_1, \dots, a_n)$.

For $n \geq 1$, an *n -ary relation* (or *predicate*) is a subset of $R \subseteq A^n$. We write $R\vec{a}$ for $\vec{a} \in R$ and $\neg R\vec{a}$ for $\vec{a} \notin R$. For $n \geq 0$, an *n -ary operation* is a function $f : A^n \rightarrow A$. A 0-ary operation is a function $f : \{\emptyset\} \rightarrow A$, and is uniquely determined by $c = f(\emptyset) \in A$. In this case, we identify f with c . A 0-ary operation is also called a *constant*.

An *extralogical signature* is a set L of symbols. Each symbol in L is called an *extralogical symbol*, and has both a *type* and an *arity*. The possible types are *relation symbols* and *function (or operation) symbols*. The arity is a nonnegative integer. Relation symbols may have an arity $n \geq 1$. Function symbols may have an arity $n \geq 0$. A 0-ary function symbol is also called a *constant symbol*.

When referring to symbols in L , we will adopt the convention that, unless otherwise stated, c will denote a constant symbol, r a relation symbol, and f a function symbol with arity $n \geq 1$.

Let L be an extralogical signature and let A be a set. For each symbol $s \in L$, let s^ω be a relation such that

- (i) if s is an n -ary relation symbol, then $s^\omega \subseteq A^n$ is an n -ary relation on A , and
- (ii) if s is an n -ary function symbol, then $s^\omega : A^n \rightarrow A$ is an n -ary function.

Let $L^\omega = \{s^\omega \mid s \in L\}$ and $\omega = (A, L^\omega)$. Then ω is an *L -structure*. The set A is called the *domain* of ω . The structure ω is called *finite* or *infinite* if A is finite or infinite, respectively.

Let L be an extralogical signature, $\nu = (B, L^\nu)$ an L -structure, and $A \subseteq B$. Suppose that for all functions $f^\nu \in L^\nu$, the subset A is *closed* under f^ν , meaning that $\vec{a} \in A^n$ implies $f^\nu \vec{a} \in A$. Let $f^\omega = f^\nu|_{A^n}$, $r^\omega = r^\nu \cap A^n$, and

$L^\omega = \{\mathfrak{s}^\omega \mid \mathfrak{s} \in L\}$. The $\omega = (A, L^\omega)$ is a structure. We call ω a *substructure* of ν , and write $\omega \subseteq \nu$.

Let L be an extralogical signature, $\omega = (A, L^\omega)$ an L -structure, and $L_0 \subseteq L$. Define $L^{\omega_0} = \{\mathfrak{s}^\omega \mid \mathfrak{s} \in L_0\}$ and $\omega_0 = (A, L^{\omega_0})$. Then ω_0 is an L_0 -structure. We call ω_0 the L_0 -*reduct* of ω , and we call ω an L -*expansion* of ω_0 .

2.4.1 Structure homomorphisms

Let $\omega = (A, L^\omega)$ and $\nu = (B, L^\nu)$ be L -structures and let $g : A \rightarrow B$. We will abuse notation and also write $g : \omega \rightarrow \nu$. For $\vec{a} \in A^n$, we write $g\vec{a}$ for $\langle ga_1, \dots, ga_n \rangle$. Assume that

- (i) $gf^\omega \vec{a} = f^\nu g\vec{a}$ for all function symbols $f \in L$,
- (ii) $gc^\omega = c^\nu$ for all constant symbols $c \in L$, and
- (iii) $r^\omega \vec{a}$ implies $r^\nu g\vec{a}$ for all relation symbols $r \in L$.

Then g is a *homomorphism*. A *strong homomorphism* is a homomorphism such that $r^\nu g\vec{a}$ implies $r^\omega \vec{b}$ for some $\vec{b} \in g^{-1}g\vec{a}$. An *embedding* is an injective strong homomorphism, an *isomorphism* is a bijective strong homomorphism, and an *automorphism* is an isomorphism from ω to ω . If g is an isomorphism, then $r^\omega \vec{a}$ if and only if $r^\nu g\vec{a}$. We say that ω and ν are *isomorphic*, written $\omega \simeq \nu$, if there is an isomorphism $g : \omega \rightarrow \nu$.

Let $\omega = (A, L^\omega)$ be a structure, let B be a set with the same cardinality as A , and let $g : A \rightarrow B$ be any bijection. Define the L -structure ν with domain B by

- (i)' $f^\nu \vec{b} = gf^\omega g^{-1} \vec{b}$ for all function symbols $f \in L$,
- (ii)' $c^\nu = gc^\omega$ for all constant symbols $c \in L$, and
- (iii)' $r^\nu \vec{b}$ if and only if $r^\omega g^{-1} \vec{b}$ for all relation symbols $r \in L$.

Then g is an isomorphism from ω to ν , and we call ν the *isomorphic image* of ω under g .

2.4.2 The standard structure of arithmetic

(S:std-arith) Consider the extralogical signature $L = \{\underline{0}, \mathfrak{S}, +, \cdot, <\}$, where $\underline{0}$ is a constant symbol, \mathfrak{S} is a unary function symbol, $+$ and \cdot are binary function symbols, and $<$ is a binary relation symbol.

We can define the structure $\mathcal{N} = (\mathbb{N}_0, L^\mathcal{N})$ as follows. Let $\underline{0}^\mathcal{N} = 0$, and let $\mathfrak{S}^\mathcal{N}$ be the successor function, $n \mapsto n + 1$. Define $+\mathcal{N}$ and $\cdot^\mathcal{N}$ to be ordinary addition and multiplication in \mathbb{N}_0 , and $<^\mathcal{N}$ to be the ordinary less-than relation on \mathbb{N}_0 . For $+$, \cdot , and $<$, we will need to rely on context to distinguish between the extralogical symbols and their ordinary meanings. To make matters worse, we will also sometimes use \mathfrak{S} to denote the function $n \mapsto n + 1$, so that context is also required to determine whether \mathfrak{S} denotes a symbol or a function.

The structure $\mathcal{N} = (\mathbb{N}_0, 0, \mathbf{S}, +, \cdot, <)$ is called the *standard structure of arithmetic*. Sometimes, we will use this phrase to refer to the same structure, but with $<$ omitted.

Now consider a different structure $\mathcal{N}' = (\mathbb{N}'_0, L^{\mathcal{N}'})$, defined as follows. Let $0^{\mathcal{N}'} = \emptyset$. Let $\mathbf{S}^{\mathcal{N}'} = s$, where $s(\alpha) = \alpha \cup \{\alpha\}$ is the successor function on ordinals. Define $+^{\mathcal{N}'}$ and $\cdot^{\mathcal{N}'}$ to be ordinal addition and ordinal multiplication, and let $<^{\mathcal{N}'} = \in$. It can be shown that the function $n \mapsto n'$, defined in Section 2.1.1, is an isomorphism from \mathcal{N} to \mathcal{N}' . We may therefore sometimes identify \mathcal{N} and \mathcal{N}' , thinking of the natural numbers as being identical to the finite ordinals.

2.4.3 Factor structures

Let $\omega = (A, L^\omega)$ be a structure and let \approx be an equivalence relation on A . We write $\vec{a} \approx \vec{b}$ to mean that $a_i \approx b_i$ for all i . Suppose that for all function symbols $f \in L$, we have $\vec{a} \approx \vec{b}$ implies $f^\omega \vec{a} = f^\omega \vec{b}$. Then \approx is a *congruence (relation) in ω* .

Let ω and ν be structures and $g : \omega \rightarrow \nu$ a homomorphism. Define $\approx_g \subseteq A^2$ by $a \approx_g b$ if and only if $ga = gb$. Then \approx_g is a congruence in ω called the *kernel of g* . Let $A' = A/\approx$ be the set of equivalence classes under \approx . Let a/\approx denote the equivalence class of a and write $\vec{a}/\approx = \langle a_1/\approx, \dots, a_n/\approx \rangle$.

If $f \in L$ is a function symbol, define $f^{\omega'} : (A')^n \rightarrow A'$ by $f^{\omega'}(\vec{a}/\approx) = (f^\omega \vec{a})/\approx$. If $r \in L$ is a relation symbol, define $r^{\omega'} \subseteq (A')^n$ by $r^{\omega'}(\vec{a}/\approx)$ if and only if $r^\omega \vec{b}$ for some $\vec{b} \approx \vec{a}$. It can be shown that both $f^{\omega'}$ and $r^{\omega'}$ are well-defined. We also define $c^{\omega'} = c^\omega/\approx$. Then $\omega' = (A', L^{\omega'})$ is a structure, denoted by ω/\approx , and called the *factor structure of ω modulo \approx* .

Let \approx be a congruence in an L -structure ω . According to the homomorphism theorem (see, for example, [28, Section 2.1]), the map $a \mapsto a/\approx$ is a strong homomorphism from ω onto ω/\approx , which we call the *canonical homomorphism*.

Conversely, let ω and ν be L -structures and $g : \omega \rightarrow \nu$ a surjective strong homomorphism. Let \approx be the kernel of g . Let k be the canonical homomorphism from $\omega \rightarrow \omega/\approx$ and let ι denote the map $a/\approx \mapsto ga$. Also according to the homomorphism theorem, the map ι is an isomorphism from ω/\approx to ν , and $g = \iota \circ k$.

2.4.4 Direct products of structures

Let L be an extralogical signature and let $\langle \omega_i \mid i \in I \rangle$ be an indexed collection of L -structures. We let $B = \prod_{i \in I} A_i$ and adopt the notation $a = \langle a_i \mid i \in I \rangle \in B$, $\vec{a} = \langle a^1, \dots, a^n \rangle \in B^n$, and $\vec{a}_i = \langle a_i^1, \dots, a_i^n \rangle \in A_i^n$. For symbols in L , we define relations and operations on B by $r^\nu \vec{a}$ if and only if $r^{\omega_i} \vec{a}_i$ for all $i \in I$, $f^\nu \vec{a} = \langle f^{\omega_i} \vec{a}_i \mid i \in I \rangle$, and $c^\nu = \langle c^{\omega_i} \mid i \in I \rangle$. Then $\nu = (B, L^\nu)$ is a structure called the *direct product* of $\langle \omega_i \mid i \in I \rangle$, and denoted by $\prod_{i \in I} \omega_i$.

If $\omega_i = \omega$ for all $i \in I$, then ν is called the *direct power* and is denoted by ω^I . If $I = \{1, \dots, n\}$, then we denote $\prod_{i \in I} \omega_i$ by $\omega_1 \times \dots \times \omega_n$, and we denote ω^I by ω^n . Note that $a \mapsto \langle a \mid i \in I \rangle$ is an embedding from ω to ω^I .

2.5 Strings

Let A be a nonempty set, which we will call an *alphabet*. Formally, it does not matter what the elements of A are, but in this context, we will refer to the elements of A as *symbols*. The set of (*finite*) *strings* over A is the set, $S = \bigcup_{n=0}^{\infty} A^n$. A string in A^n is said to have *length* n . The unique string of length 0 is \emptyset , and we call this the *empty string*. A string of length 1 is called an *atomic* string.

Strings are written without brackets or commas, so that we write $s_1 \cdots s_n$, rather than $\langle s_1, \dots, s_n \rangle$. Let $\xi = s_1 \cdots s_n$ and $\eta = s_{n+1} \cdots s_{n+m}$ be strings of lengths n and m , respectively. The *concatenation* of ξ and η , denoted by $\xi\eta$, is the string $s_1 \cdots s_{n+m}$.

If $\xi = \xi_1\eta\xi_2$, where $\eta \neq \emptyset$, then η is called a *segment* (or *substring*) of ξ . If $\eta \neq \xi$, then η is a *proper* segment (or substring) of ξ . If $\xi_1 = \emptyset$, then η is an *initial* segment (or substring) of ξ . If $\xi_2 = \emptyset$, then η is a *terminal* segment (or substring) of ξ .

Chapter 3

Propositional Calculus

(Ch:prop-calc) In this chapter, we develop a calculus of inductive inference for sentences in a formal logical language. The language we focus on, denoted by \mathcal{F} , is propositional. The sentences (or formulas) of \mathcal{F} consist of primitive propositional variables connected by negation and conjunction. Using negation and conjunction, we can define other logical connectives, such as disjunction and material implication. Our language \mathcal{F} , however, is not the usual propositional language one finds in basic logic textbooks. Our language is infinitary, in the sense that we allow countable conjunctions. That is, if φ_n is a formula of \mathcal{F} for every n , then so is $\bigwedge_1^\infty \varphi_n$.

The calculus of deductive inference in \mathcal{F} is well understood. The study of infinitary languages dates back to the papers of Scott and Tarski [30, 32] and to the dissertation [17] and book [16] of Carol Karp. Deductive calculus in \mathcal{F} is represented by a derivability relation, \vdash . If φ is a formula of \mathcal{F} and X is a set of formulas, then $X \vdash \varphi$ denotes the fact that φ can be derived from X using the rules of deductive inference. In Section 3.1, we define \mathcal{F} and \vdash , and establish important facts about them. We then present the notion of a deductive theory, which is a set of formulas that is closed under deductive inference. An important point, noted in Remark 3.1.21, is that one can define deductive theories without any reference to derivability, and then define derivability in terms of theories. This is the approach we take in the development of our inductive calculus.

An inductive statement is a triple, (X, φ, p) , where X is a set of formulas, φ is a formula, and $p \in [0, 1]$. Intuitively, an inductive statement can be thought of as asserting that φ is partially entailed by X with degree p . In an inductive statement, X is called the antecedent, φ is called the consequent, and p is called the probability. We typically use a letter such as P to denote sets of inductive statements. We write $P(\varphi \mid X) = p$ to mean that $(X, \varphi, p) \in P$ and refer to p as the probability of φ given X . The ultimate goal of this chapter is to develop a calculus by which we can take a given set of inductive statements P , and infer a new inductive statement (X, φ, p) . When such an inference is possible, we will write $P \vdash (X, \varphi, p)$, thereby extending the use of the turnstile symbol \vdash from formulas to inductive statements.

As mentioned above, we take an indirect route to defining \vdash . We first define an inductive theory. Intuitively, if P is an inductive theory, and if it is possible to infer (X, φ, p) from P , then (X, φ, p) is already an element of P . The bulk of this chapter is devoted to defining the notion of an inductive theory. Once this is done, we define \vdash in terms of inductive theories.

In Sections 3.2 and 3.3, we present our rules of inductive inference. There are nine of them altogether, which we denote by (R1)–(R9). Rules (R1)–(R4) describe the connection between inductive and deductive inference. Rules (R5)–(R7) are the core rules of inductive inference: the addition, multiplication, and continuity rules. Rule (R8) says that we can use uniqueness to make inferences: if there is a unique way to assign a probability without violating (R1)–(R7), then we may infer that probability. Rule (R9) says that we may freely add formulas of probability one to our antecedents.

To define an inductive theory, we must first define what it means for a set of inductive statements to be closed under the rules of inference. We do this in tiers. An admissible set is one that is closed under (R1), an entire set is closed under (R1)–(R7), a semi-closed set is closed under (R1)–(R8), and a closed set is closed under (R1)–(R9).

In Section 3.2, we define admissible and entire sets. We then prove several theorems about entire sets, such as inclusion-exclusion, Bayes' theorem, and countable additivity. Section 3.3 begins with the definition of semi-closed and closed sets. We then turn our attention to a notion that has no analogue in the deductive calculus. The set of inductive statements is much larger than the set of formulas. As such, it is possible to have inductive statements that are so far apart, in a certain sense, that they can never be related to one another via the rules of inductive inference. That is, no chain of inductive reasoning could possibly include both statements. Such statements would be trivially incapable of producing a contradiction. Yet they are also incapable of meaningfully contributing to a common argument. Because of this, a closed set could potentially contain components that bear no logical connection to one another. In Section 3.3, we make this notion of connectivity precise. We then define an inductive theory to be a closed set that is also connected. Section 3.3 concludes by using the definition of an inductive theory to define the inductive derivability relation, $P \vdash (X, \varphi, p)$.

The theorems in Section 3.3 are presented without proof, so that the reader can see a complete overview of the development. Their proofs are presented in Sections 3.4 and 3.5. In Section 3.4, we prove that inductive theories exist and are well-defined. The proof is primarily constructive, showing how to build up an inductive theory from more basic elements. We then say that a set of inductive statements is consistent if it can be extended to an inductive theory. In Section 3.5, we prove that every consistent set can be uniquely extended to an inductive theory.

Section 3.5 concludes with an important generalization of inductive derivability. So far we have only discussed a process of inference whereby we take a set of inductive statements, P , and use them to derive a new inductive statement, (X, φ, p) . In an inferential argument such as this, the inductive

statements in P are our hypotheses. Since every inductive statement has the form $P(\varphi \mid X) = p$, it follows that every one of our hypotheses must have this form as well. In Section 3.5, we allow for a broader class of hypotheses. By introducing what we call inductive conditions, we are able to use hypotheses such as $P(\varphi \mid X) > p$ or $P(\varphi \wedge \psi \mid X, \psi) = P(\varphi \mid X)$.

3.1 Formulas and deductive inference

$\langle S: \text{formulas} \rangle$

3.1.1 Propositional formulas

Let PV be a nonempty set whose elements we call *propositional variables*. We define an alphabet, $A = PV \cup \{\neg, \wedge\}$. We will define the set of formulas so that a formula is a finite tuple, where each element in the tuple is either a symbol from our alphabet, a formula, or a countable set of formulas.

Let $S_0 = \{\langle p \rangle \mid p \in PV\}$. For an ordinal $\alpha < \omega_1$, let

$$S'_\alpha = S_\alpha \cup \{\langle \neg, \varphi \rangle \mid \varphi \in S_\alpha\}.$$

When writing tuples such as these, we will typically omit the commas and angled brackets, so that, for instance, $\neg\varphi = \langle \neg, \varphi \rangle$. In particular, for $\mathbf{r} \in PV$, we identify $\langle \mathbf{r} \rangle$ and \mathbf{r} so that we may write $S_0 = PV$.

We then define

$$S_{\alpha+1} = S'_\alpha \cup \{\langle \wedge, \Phi \rangle \mid \Phi \subseteq S'_\alpha \text{ is nonempty and countable}\}.$$

Here, countable means finite or countably infinite. As above, we will typically write $\wedge\Phi$ as shorthand for ordered pairs of this type.

In the case that α is a limit ordinal, we define $S_\alpha = \bigcup_{\xi < \alpha} S_\xi$. Finally, we define $\mathcal{F} = \mathcal{F}_{\omega_1} = \bigcup_{\alpha < \omega_1} S_\alpha$. Note that $S_\alpha \subseteq S_\beta$ whenever $\alpha < \beta$. An element $\varphi \in \mathcal{F}$ is called a (*propositional*) *formula* or *sentence*. A formula may also be called a *deductive statement*, in contrast to inductive statements, to be defined later. The set \mathcal{F} depends on the choice of PV . We will rarely need to emphasize this fact, but when we do, we will write $\mathcal{F}(PV)$ instead of \mathcal{F} .

Let \mathcal{F}_{fin} denote the smallest subset of \mathcal{F} that satisfies

- (i) $PV \subseteq \mathcal{F}_{\text{fin}}$,
- (ii) if $\varphi \in \mathcal{F}_{\text{fin}}$, then $\neg\varphi \in \mathcal{F}_{\text{fin}}$, and
- (iii) if $\Phi \subseteq \mathcal{F}_{\text{fin}}$ is nonempty and finite, then $\wedge\Phi \in \mathcal{F}_{\text{fin}}$.

Formulas in \mathcal{F}_{fin} are said to be *finitary*. The set \mathcal{F}_{fin} is, in fact, the set of formulas used in finitary propositional logic. Or rather, it is one of several equivalent definitions of the finitary propositional language. The reader can consult any introductory text on mathematical logic for the basic properties of \mathcal{F}_{fin} and its corresponding syntax and semantics. When necessary, we will cite [28] for this purpose.

(T:prop-form-ind) **Theorem 3.1.1 (The principle of formula induction).** *The set of formulas, \mathcal{F} , is the smallest set that satisfies the following:*

- (i) $PV \subseteq \mathcal{F}$,
- (ii) if $\varphi \in \mathcal{F}$, then $\neg\varphi \in \mathcal{F}$, and
- (iii) if $\Phi \subseteq \mathcal{F}$ is nonempty and countable, then $\bigwedge \Phi \in \mathcal{F}$.

Proof. Property (i) follows since $PV = S_0 \subseteq \mathcal{F}$. Suppose $\varphi \in \mathcal{F}$. Then there exists $\alpha < \omega_1$ such that $\varphi \in S_\alpha$. Thus, $\neg\varphi \in S'_\alpha \subseteq \mathcal{F}$, proving property (ii). Finally, suppose $\Phi \subseteq \mathcal{F}$ is nonempty and countable. Enumerate Φ as $\Phi = \{\varphi_n\}_{n < \alpha}$, where $0 < \alpha \leq \omega$. For each $n < \alpha$, choose $\alpha_n < \omega_1$ such that $\varphi_n \in S_{\alpha_n}$. Choose $\beta < \omega_1$ such that $\alpha_n \leq \beta$ for all $n < \alpha$. It follows that Φ is a nonempty, countable subset of $S_\beta \subseteq S'_\beta$, and so $\bigwedge \Phi \in S_{\beta+1} \subseteq \mathcal{F}$, proving (iii).

Now let S be any set that satisfies these three properties. Using transfinite induction, it is easy to verify that $S_\alpha \subseteq S$ for all ordinals α . Thus, $\mathcal{F} = \bigcup_{\alpha < \omega_1} S_\alpha \subseteq S$. \square

We may sometimes write $\bigwedge_{\varphi \in \Phi} \varphi$ for $\bigwedge \Phi$. If $\Phi = \{\varphi_n\}_{n \in \mathbb{N}}$, we may write $\bigwedge_{n=1}^{\infty} \varphi_n$ for $\bigwedge \Phi$. If $\Phi \subseteq \mathcal{F}$, then we use the notation $\neg\Phi = \{\neg\varphi : \varphi \in \Phi\}$. As shorthand, we define

$$\begin{aligned} \bigvee \Phi &= \neg \bigwedge \neg \Phi \text{ for } \Phi \text{ nonempty and countable,} \\ (\varphi \wedge \psi) &= \bigwedge \{\varphi, \psi\}, \\ (\varphi \vee \psi) &= \bigvee \{\varphi, \psi\}, \\ (\varphi \rightarrow \psi) &= (\neg\varphi \vee \psi), \text{ and} \\ (\varphi \leftrightarrow \psi) &= (\varphi \rightarrow \psi) \wedge (\psi \rightarrow \varphi). \end{aligned}$$

We fix an arbitrary $\mathbf{r}_0 \in PV$ and define $\perp = (\mathbf{r}_0 \wedge \neg\mathbf{r}_0)$ and $\top = \neg\perp$. The symbols, \perp and \top , are called *falsum* and *verum*, respectively. We adopt the convention that $\bigwedge \emptyset = \top$ and $\bigvee \emptyset = \perp$.

As another form of shorthand, we may also sometimes omit outer parentheses in formulas, and occasionally inner parentheses with the understanding that \rightarrow associates right to left, other symbols associate left to right, and that formulas obey the order of operations, $\neg, \wedge, \vee, \rightarrow, \leftrightarrow$.

Finally, we may occasionally use the notation φ^x , where x is an element of the Boolean algebra $\mathbf{B} = \{0, 1\}$, to mean $\varphi^1 = \varphi$ and $\varphi^0 = \neg\varphi$.

Given $\varphi \in \mathcal{F}$, we define the set of *subformulas* of φ , denoted by $\text{Sf } \varphi$, by formula recursion. Namely, $\text{Sf } \mathbf{r} = \{\mathbf{r}\}$ for $\mathbf{r} \in PV$, $\text{Sf } \neg\varphi = \{\neg\varphi\} \cup \text{Sf } \varphi$, and $\text{Sf } \bigwedge \Phi = \{\bigwedge \Phi\} \cup \bigcup_{\varphi \in \Phi} \text{Sf } \varphi$. It follows by formula induction that $\text{Sf } \varphi$ is countable for every $\varphi \in \mathcal{F}$. Also, the set of propositional variables that appear in a formula φ is simply $PV \cap \text{Sf } \varphi$. In particular, each formula φ makes use of only countably many propositional variables.

Remark 3.1.2. The construction presented here is analogous to the one suggested in [18] for formulas in infinitary predicate logic. An alternative

construction builds formulas out of countably long sequences of symbols in an alphabet. In this case, one must deal precisely with the notion of concatenation for such strings, such as the concatenation of countably many strings, each of which may itself be countably long. All of this is detailed in [16].

3.1.2 A calculus of natural deduction

$\langle \text{S:prop-nat-ded} \rangle$ Given a relation \vdash from $\mathfrak{P}\mathcal{F}$ to \mathcal{F} , we write $X \vdash Y$ to mean $X \vdash \varphi$ for all $\varphi \in Y$. A comma-separated list on either side of the turnstile, \vdash , refers to a union, and isolated formulas refer to the singleton set that contains them. For example, $X, Y, \varphi \vdash \psi, \zeta$ means $X \cup Y \cup \{\varphi\} \vdash \{\psi, \zeta\}$, which means that $X \cup Y \cup \{\varphi\} \vdash \psi$ and $X \cup Y \cup \{\varphi\} \vdash \zeta$. Also, $\vdash \varphi$ is shorthand for $\emptyset \vdash \varphi$.

We wish to define a relation \vdash from $\mathfrak{P}\mathcal{F}$ to \mathcal{F} such that $X \vdash \varphi$ captures what it means to say that φ can be logically deduced from the formulas in X .

$\langle \text{D:derivability} \rangle$ **Definition 3.1.3.** The *derivability relation* is the smallest relation \vdash from $\mathfrak{P}\mathcal{F}$ to \mathcal{F} such that, for all $\varphi, \psi \in \mathcal{F}$ and all countable $\Phi \subseteq \mathcal{F}$, the following conditions hold:

- (i) $\varphi \vdash \varphi$,
- (ii) if $X \vdash \varphi$ and $X \subseteq X'$, then $X' \vdash \varphi$,
- (iii) if $X \vdash \bigwedge \Phi$, then $X \vdash \theta$ for all $\theta \in \Phi$,
- (iv) if $X \vdash \theta$ for all $\theta \in \Phi$, then $X \vdash \bigwedge \Phi$,
- (v) if $X \vdash \varphi$ and $X \vdash \neg\varphi$, then $X \vdash \psi$, and
- (vi) if $X, \varphi \vdash \psi$ and $X, \neg\varphi \vdash \psi$, then $X \vdash \psi$.

When $X \vdash \varphi$, we say that φ is (*deductively*) *derivable* from X , or that X *proves* φ . If $X \not\vdash \varphi$ and $X \not\vdash \neg\varphi$, then φ is *undetermined by* (or *deductively independent of*) X . We say that φ is *determined by* X if φ is not undetermined by X .

Remark 3.1.4. Note that the intersection of any family of relations that satisfy (i)–(vi) also satisfies (i)–(vi). Also note that \mathcal{F} itself satisfies (i)–(vi). Thus, the derivability relation is well-defined and is equal to the intersection of all subsets of $\mathfrak{P}\mathcal{F} \times \mathcal{F}$ satisfying (i)–(vi).

Remark 3.1.5. An alternative but equivalent definition of the derivability relation involves the notion of a derivation. A derivation of φ from X is a countable sequence $\langle (X_\beta, \varphi_\beta) : \beta \leq \alpha \rangle$, where α is a countable ordinal, $(X_\alpha, \varphi_\alpha) = (X, \varphi)$, and for each $\beta \leq \alpha$, the term (X_β, φ_β) is obtained from $\langle (X_\xi, \varphi_\xi) : \xi < \beta \rangle$ by an application of one of the six rules in Definition 3.1.3. In this case, $X \vdash \varphi$ if and only if there exists a derivation of φ from X .

Definition 3.1.6. The *finitary derivability relation* is the smallest relation \vdash_{fin} from $\mathfrak{P}\mathcal{F}_{\text{fin}}$ to \mathcal{F}_{fin} such that, for all $\varphi, \psi \in \mathcal{F}_{\text{fin}}$ and all finite $\Phi \subseteq \mathcal{F}$, conditions (i)–(vi) from Definition 3.1.3 hold.

$\langle R:\text{finitary-vs-infinitary} \rangle$ **Remark 3.1.7.** The finitary derivability relation is a typical natural-deduction calculus for finitary propositional logic. Clearly, $\vdash_{\text{fin}} \subseteq \vdash$. As we will see in Proposition 4.1.15, if $X \subseteq \mathcal{F}_{\text{fin}}$, $\varphi \in \mathcal{F}_{\text{fin}}$, and $X \vdash \varphi$, then $X \vdash_{\text{fin}} \varphi$. In other words, when restricted to finitary formulas, our infinitary calculus cannot produce any new inferences beyond those already available with the finitary calculus.

The proofs of the structural rules in the following proposition are exactly the same as in finitary propositional logic (see, for instance, [28, Section 1.4]).

$\langle P:\text{derivability} \rangle$ **Proposition 3.1.8.** *The derivability relation satisfies the following:*

- (a) $X \vdash (\varphi \rightarrow \psi)$ if and only if $X, \varphi \vdash \psi$,
- (b) if $X \vdash \varphi$ and $X, \varphi \vdash \psi$, then $X \vdash \psi$,
- (c) if $X, \neg\varphi \vdash \varphi$, then $X \vdash \varphi$, and
- (d) if $X, \varphi \vdash \neg\varphi$, then $X \vdash \neg\varphi$.

In the remainder of this section, unless otherwise specified, lowercase Roman numerals refer to Definition 3.1.3 and letters refer to Proposition 3.1.8.

$\langle L:\text{conj-subset} \rangle$ **Lemma 3.1.9.** *Let $\Phi, \Phi' \subseteq \mathcal{F}$ be countable and $\psi \in \mathcal{F}$. If $\Phi \subseteq \Phi'$ and $\bigwedge \Phi \vdash \psi$, then $\bigwedge \Phi' \vdash \psi$.*

Proof. Suppose $\Phi \subseteq \Phi'$ and $\bigwedge \Phi \vdash \psi$. By (i) and (iii), we have $\bigwedge \Phi' \vdash \theta$ for all $\theta \in \Phi$. Thus, by (iv), we have $\bigwedge \Phi' \vdash \bigwedge \Phi$. The result now follows from (b). \square

$\langle T:\text{sig-cpctness} \rangle$ **Theorem 3.1.10 (σ -compactness).** *Let $X \subseteq \mathcal{F}$ and $\varphi \in \mathcal{F}$. Then $X \vdash \varphi$ if and only if there exists a countable subset $X_0 \subseteq X$ such that $X_0 \vdash \varphi$.*

Proof. We will actually prove that the following are equivalent:

$$X \vdash \varphi, \tag{3.1.1} \boxed{\text{sig-cpctness-1}}$$

$$\text{there exists countable } X_0 \subseteq X \text{ such that } \bigwedge X_0 \vdash \varphi, \text{ and} \tag{3.1.2} \boxed{\text{sig-cpctness-2}}$$

$$\text{there exists countable } X_0 \subseteq X \text{ such that } X_0 \vdash \varphi. \tag{3.1.3} \boxed{\text{sig-cpctness-3}}$$

By (i), (ii), (iv), and (b), we have that (3.1.2) implies (3.1.3), and by (ii), we have (3.1.3) implies (3.1.1).

To prove that (3.1.1) implies (3.1.2), we define \vdash' so that $X \vdash' \varphi$ if and only if $X \vdash \varphi$ and (3.1.2) holds. The proof will be complete once we show that (i)–(vi) still hold when \vdash is replaced by \vdash' .

Clearly, (i)–(iii) hold for \vdash' . To see that (iv)–(vi) hold for \vdash' , use Lemma 3.1.9 and the fact that a countable union of countable sets is countable. \square

$\langle P:\text{conj-equiv} \rangle$ **Proposition 3.1.11.** *Let $\Phi \subseteq \mathcal{F}$ be countable and $\psi \in \mathcal{F}$. Then $\Phi \vdash \psi$ if and only if $\bigwedge \Phi \vdash \psi$.*

Proof. By (i), (ii), and (iv), we have $\Phi \vdash \bigwedge \Phi$. Thus, by (b), it follows that $\bigwedge \Phi \vdash \psi$ implies $\Phi \vdash \psi$.

Now suppose $\Phi \vdash \psi$. By (3.1.2), there exists countable $\Phi' \subseteq \Phi$ such that $\bigwedge \Phi' \vdash \psi$. By Lemma 3.1.9, we have $\bigwedge \Phi \vdash \psi$. \square

$\langle \text{P:set-mono} \rangle$ **Proposition 3.1.12.** *If $X \vdash Y$ and $Y \vdash \varphi$, then $X \vdash \varphi$.*

Proof. Suppose $X \vdash Y$ and $Y \vdash \varphi$. By σ -compactness, there exists countable $Y_0 \subseteq Y$ such that $Y_0 \vdash \varphi$. By Proposition 3.1.11, we have $\bigwedge Y_0 \vdash \varphi$. By (iv), we have $X \vdash \bigwedge Y_0$. Hence, by (b), we have $X \vdash \varphi$. \square

A set $X \subseteq \mathcal{F}$ is *inconsistent* if $X \vdash \perp$ for all $\varphi \in \mathcal{F}$; it is otherwise *consistent*. If X is inconsistent, then $X \vdash \perp$. Conversely, by the definition of \perp , and by (iii) and (v), we have that $X \vdash \perp$ implies X is inconsistent. Thus, X is inconsistent if and only if $X \vdash \perp$. The derivability relation can, in fact, be characterized in terms of consistency.

$\langle \text{T:deduc-con} \rangle$ **Theorem 3.1.13.** *Suppose $X \subseteq \mathcal{F}$ and $\varphi \in \mathcal{F}$. Then $X \vdash \varphi$ if and only if $X, \neg\varphi \vdash \perp$, and $X \vdash \neg\varphi$ if and only if $X, \varphi \vdash \perp$.*

Proof. Suppose $X \vdash \varphi$. By (i) and (ii), we have $X, \neg\varphi \vdash \varphi$ and $X, \neg\varphi \vdash \neg\varphi$. By (v), this implies $X, \neg\varphi \vdash \perp$. Conversely, suppose $X, \neg\varphi \vdash \perp$. Then $X \cup \{\neg\varphi\}$ is inconsistent, so that $X, \neg\varphi \vdash \varphi$. By (c), we have $X \vdash \varphi$. The proof of the second biconditional is analogous. \square

Since $\top = \neg\perp$, the preceding theorem shows that $X \vdash \top$ if and only if $X, \perp \vdash \perp$. Hence, by (i) and (ii), we have $X \vdash \top$ for all $X \subseteq \mathcal{F}$, which by (ii) is equivalent to $\vdash \top$.

A formula φ is a *tautology* if $\vdash \varphi$; it is a *contradiction* if $\{\varphi\}$ is inconsistent. By the preceding proposition, we see that φ is a tautology if and only if $\neg\varphi$ is a contradiction, and vice versa. The set of tautologies is denoted by Taut , or $\mathit{Taut}_{\mathcal{F}}$.

$\langle \text{P:sig-cpctness} \rangle$ **Proposition 3.1.14.** *Let $X \subseteq \mathcal{F}$ and $\varphi \in \mathcal{F}$. Then $X \vdash \varphi$ if and only if there exists a countable $X_0 \subseteq X$ such that $\bigwedge X_0 \rightarrow \varphi \in \mathit{Taut}$.*

Proof. By σ -compactness, we have $X \vdash \varphi$ if and only if there exists countable $X_0 \subseteq X$ such that $X_0 \vdash \varphi$. And by Proposition 3.1.11 and (a), we have $X_0 \vdash \varphi$ if and only if $\bigwedge X_0 \rightarrow \varphi \in \mathit{Taut}$. \square

3.1.3 A Hilbert-type calculus

$\langle \text{S:Hilbert-calc} \rangle$ Let $\Lambda = \Lambda_{\mathcal{F}}$ be the smallest subset of \mathcal{F} such that if $\varphi, \psi, \zeta \in \mathcal{F}$ and $\Phi \subseteq \mathcal{F}$ is countable with $\varphi \in \Phi$, then the following formulas are in Λ :

$$(\Lambda 1) \quad (\varphi \rightarrow \psi \rightarrow \zeta) \rightarrow (\varphi \rightarrow \psi) \rightarrow \varphi \rightarrow \zeta$$

$$(\Lambda 2) \quad (\varphi \rightarrow \neg\psi) \rightarrow \psi \rightarrow \neg\varphi$$

$$(\Lambda 3) \quad \bigwedge \Phi \rightarrow \varphi$$

The formulas in Λ are called *axioms*.

We define a *proof of $\varphi \in \mathcal{F}$ from $X \subseteq \mathcal{F}$* as an $(\alpha + 1)$ -sequence of formulas, $\langle \varphi_\beta \mid \beta \leq \alpha \rangle$, where α is a countable ordinal, $\varphi_\alpha = \varphi$, and for each $\beta \leq \alpha$, either $\varphi_\beta \in X \cup \Lambda$, or there exist $i, j < \beta$ such that $\varphi_i = (\varphi_j \rightarrow \varphi_\beta)$, or there exists nonempty, countable $\Phi \subseteq \{\varphi_\xi \mid \xi < \beta\}$ such that $\varphi_\beta = \bigwedge \Phi$. Note that if $\langle \varphi_\beta \mid \beta \leq \alpha \rangle$ is a proof of φ_α from X , then for any $\beta < \alpha$, it follows that $\langle \varphi_\xi \mid \xi \leq \beta \rangle$ is a proof of φ_β from X . For $\varphi \in \mathcal{F}$ and $X \subseteq \mathcal{F}$, define $X \vdash \varphi$ to mean there is a proof of φ from X .

Let Λ_{fin} be the smallest subset of \mathcal{F}_{fin} such that if $\varphi, \psi, \zeta \in \mathcal{F}_{\text{fin}}$ and $\Phi \subseteq \mathcal{F}_{\text{fin}}$ is finite with $\varphi \in \Phi$, then $(\Lambda 1)$ – $(\Lambda 3)$ are in Λ_{fin} . A *finitary proof of $\varphi \in \mathcal{F}_{\text{fin}}$ from $X \subseteq \mathcal{F}_{\text{fin}}$* is a finite sequence of formulas, $\langle \varphi_k \mid k \leq n \rangle$, where $\varphi_n = \varphi$, and for each $k \leq n$, either $\varphi_k \in X \cup \Lambda_{\text{fin}}$, or there exist $i, j < k$ such that $\varphi_i = (\varphi_j \rightarrow \varphi_k)$, or there exists nonempty, finite $\Phi \subseteq \{\varphi_\ell \mid \ell < k\}$ such that $\varphi_k = \bigwedge \Phi$. For $\varphi \in \mathcal{F}_{\text{fin}}$ and $X \subseteq \mathcal{F}_{\text{fin}}$, define $X \vdash_{\text{fin}} \varphi$ to mean there is a finitary proof of φ from X .

A finitary proof is the classical notion of proof. It is finitely long, and each sentence in it has finite length. An infinitary proof, on the other hand, can be infinitely long. And individual sentences in such a proof can themselves have infinite length.

(R:fin-Hilbert) **Remark 3.1.15.** The relation \vdash_{fin} is a typical Hilbert-style calculus for finitary propositional logic. It is well-known that $\vdash_{\text{fin}} = \vdash$. (See, for example, [28, Theorem 1.6.6].) In Theorem 3.1.17, we will see that $\sim = \vdash$. Hence, according to Remark 3.1.7, if $X \subseteq \mathcal{F}_{\text{fin}}$, $\varphi \in \mathcal{F}_{\text{fin}}$, and $X \sim \varphi$, then $X \vdash_{\text{fin}} \varphi$. In other words, if we can find an infinitary proof of φ from X , then a finitary proof necessarily exists.

(P:Hilbert-induc) **Proposition 3.1.16 (Induction principle for \sim).** *The relation \sim is the smallest relation from $\mathfrak{P}\mathcal{F}$ to \mathcal{F} such that if $X \subseteq \mathcal{F}$, $\varphi, \psi \in \mathcal{F}$, and $\Phi \subseteq \mathcal{F}$ is countable, then*

- (1) $X \sim \theta$ for all $\theta \in X \cup \Lambda$,
- (2) if $X \sim (\varphi \rightarrow \psi)$ and $X \sim \varphi$, then $X \sim \psi$, and
- (3) if $X \sim \theta$ for all $\theta \in \Phi$, then $X \sim \bigwedge \Phi$.

Proof. The fact that \sim satisfies (1)–(3) follows from the fact that a countable concatenation of proofs is again a proof. (See [16, Chapter 2] for details on infinitary concatenation.)

Let \triangleright be a relation from $\mathfrak{P}\mathcal{F}$ to \mathcal{F} satisfying (1)–(3). Let $X \subseteq \mathcal{F}$ be arbitrary. Fix a countable ordinal α , and consider the statement,

for all $\varphi \in \mathcal{F}$, if there exists a proof of φ from X with length $\alpha + 1$,
then $X \triangleright \varphi$.

We will prove this statement is true for all countable α by induction on α , and this will show that $X \sim \varphi$ implies $X \triangleright \varphi$.

If φ has a proof from X of length 1, then it must be $\langle \varphi \rangle$, implying $\varphi \in X \cup \Lambda$. Therefore, $X \triangleright \varphi$ by (1), and the statement is true for $\alpha = 0$.

Suppose the statement is true for all $\beta < \alpha$, and that φ has a proof from X of length $\alpha + 1$. Let $\langle \varphi_\beta \mid \beta \leq \alpha \rangle$ be such a proof. If $\varphi \in X \cup \Lambda$, then $X \triangleright \varphi$ by (1). Suppose there exists $i, j < \alpha$ such that $\varphi_i = (\varphi_j \rightarrow \varphi)$. Then $\langle \varphi_\beta \mid \beta \leq i \rangle$ and $\langle \varphi_\beta \mid \beta \leq j \rangle$ are proofs of φ_i and φ_j , respectively, each with length less than $\alpha + 1$. By the inductive hypothesis, $X \triangleright \varphi_i$ and $X \triangleright \varphi_j$, so that by (2), we have $X \triangleright \varphi$. Finally, suppose there exists $\Phi \subseteq \{\varphi_\beta : \beta < \alpha\}$ such that $\varphi = \bigwedge \Phi$. Each $\varphi_\beta \in \Phi$ has a proof, $\langle \varphi_\xi \mid \xi \leq \beta \rangle$, of length $\beta + 1 < \alpha + 1$, so by the inductive hypothesis, $X \triangleright \theta$ for all $\theta \in \Phi$. Thus, by (3), we have $X \triangleright \varphi$. \square

(T:Hilbert=nat) **Theorem 3.1.17.** *Let $X \subseteq \mathcal{F}$ and $\varphi \in \mathcal{F}$. Then $X \vdash \varphi$ if and only if $X \triangleright \varphi$.*

Proof. In this proof, Arabic numerals will refer to Proposition 3.1.16.

We first prove that $X \vdash \varphi$ implies $X \triangleright \varphi$. By Proposition 3.1.16, it suffices to show that (1)–(3) hold when \vdash is replaced by \triangleright .

By (iv), we have that (3) holds for \triangleright . Suppose $X \triangleright (\varphi \rightarrow \psi)$ and $X \triangleright \varphi$. By (a), we have $X, \varphi \triangleright \psi$, so by (b), it follows that $X \triangleright \psi$. Thus, (2) holds for \triangleright . By (i) and (ii), we have $X \triangleright \theta$ for all $\theta \in X$. It remains only to show that $X \triangleright \varphi$ whenever φ is an axiom. By (ii), it suffices to show that $\triangleright \varphi$ whenever φ is an axiom.

Consider first ($\Lambda 1$). Let $Y = \{\varphi \rightarrow \psi \rightarrow \zeta, \varphi \rightarrow \psi, \varphi\}$. By (i) and (ii), we have $Y \triangleright (\varphi \rightarrow \psi)$, so that (a) yields $Y, \varphi \triangleright \psi$. But $Y \cup \{\varphi\} = Y$, so $Y \triangleright \psi$. We similarly obtain $Y \triangleright (\psi \rightarrow \zeta)$, so that $Y, \psi \triangleright \zeta$. By (b), we obtain $Y \triangleright \zeta$. Repeated applications of (a) now yield $\varphi \rightarrow \psi \rightarrow \zeta, \varphi \rightarrow \psi \triangleright \varphi \rightarrow \zeta$, followed by $\varphi \rightarrow \psi \rightarrow \zeta \triangleright (\varphi \rightarrow \psi) \rightarrow \varphi \rightarrow \zeta$, followed by $\triangleright (\varphi \rightarrow \psi \rightarrow \zeta) \rightarrow (\varphi \rightarrow \psi) \rightarrow \varphi \rightarrow \zeta$.

For ($\Lambda 2$), let $X = \{\varphi \rightarrow \neg\psi, \psi\}$ and $Y = X \cup \{\varphi\}$. By (i) and (ii), we have $Y \triangleright \varphi \rightarrow \neg\psi$, so that (a) yields $Y, \varphi \triangleright \neg\psi$. But $Y \cup \{\varphi\} = Y$, so $Y \triangleright \neg\psi$. On the other hand, (i) and (ii) imply $Y \triangleright \psi$, so by (v), we have $Y \triangleright \neg\varphi$, or in other words, $X, \varphi \triangleright \neg\varphi$. By (d), we have $X \triangleright \neg\varphi$. Repeated applications of (a) yield $\varphi \rightarrow \neg\psi \triangleright \psi \rightarrow \neg\varphi$, followed by $\triangleright (\varphi \rightarrow \neg\psi) \rightarrow \psi \rightarrow \neg\varphi$.

For ($\Lambda 3$), let $\Phi \subseteq \mathcal{F}$ be countable and $\varphi \in \Phi$. By (i), we have $\bigwedge \Phi \triangleright \bigwedge \Phi$. By (iii), we obtain $\bigwedge \Phi \triangleright \varphi$. Thus, by (a), it follows that $\triangleright \bigwedge \Phi \rightarrow \varphi$.

To prove that $X \triangleright \varphi$ implies $X \vdash \varphi$, it suffices to show that (i)–(vi) hold when \triangleright is replaced by \vdash . The fact that (i), (ii), (v), and (vi) hold for \vdash follows exactly as in the finitary case (see [28, Section 1.6], for example). We obtain (iii) and (iv) from ($\Lambda 3$) and (3), respectively. \square

3.1.4 Deductive theories and logical equivalence

Definition 3.1.18. A set $T \subseteq \mathcal{F}$ is called a (*deductive*) *theory* if the following conditions hold:

- (i) $\Lambda \subseteq T$,
- (ii) if $(\varphi \rightarrow \psi) \in T$ and $\varphi \in T$, then $\psi \in T$, and

(iii) if $\Phi \subseteq T$ is countable, then $\bigwedge \Phi \in T$.

Note that the intersection of any family of theories is again a theory. Also note that \mathcal{F} itself is a theory. Hence, if $X \subseteq \mathcal{F}$, then we may define *the (deductive) theory generated by X* , denoted by $T(X)$ or T_X , as the smallest theory having X as a subset.

$\langle T:\text{theory-deduc} \rangle$ **Theorem 3.1.19.** *Let $X \subseteq \mathcal{F}$ and $\varphi \in \mathcal{F}$. Then $X \vdash \varphi$ if and only if $\varphi \in T(X)$.*

Proof. Suppose $X \vdash \varphi$. By Theorem 3.1.17, there exists a proof of φ from X . As in Proposition 3.1.16, we can use induction on the length of the proof to show that $\varphi \in T(X)$. For the converse, define $T' = \{\varphi \in \mathcal{F} : X \vdash \varphi\}$. By Theorem 3.1.17 and Proposition 3.1.16, it follows that T' is a theory with $X \subseteq T'$. Thus, $T(X) \subseteq T'$. \square

Corollary 3.1.20. *A set $T \subseteq \mathcal{F}$ is a theory if and only if it is deductively closed, meaning that $T \vdash \varphi$ implies $\varphi \in T$.*

Proof. This follows immediately by combining Theorem 3.1.19 with the fact that $T \subseteq \mathcal{F}$ is a theory if and only if $T = T(T)$. \square

$\langle R:\text{theory-first} \rangle$ **Remark 3.1.21.** Theorem 3.1.19 exhibits an alternative approach to defining derivability. One can define the theory generated by a set X , as we did above, without any reference to derivability. Then one can define derivability in terms of $T(X)$. This is the approach we will take when considering inductive inference.

If T is a theory and $S \subseteq \mathcal{F}$, then $T + S$ denotes the theory generated by $T \cup S$. For $\varphi \in \mathcal{F}$, we write $T + \varphi$ for $T + \{\varphi\}$. The smallest theory is **Taut**, the largest theory is \mathcal{F} , and every theory T satisfies **Taut** $\subseteq T \subseteq \mathcal{F}$. A theory T is inconsistent if and only if $T = \mathcal{F}$. A theory T is said to be *(deductively) complete* if it is consistent and every $\varphi \in \mathcal{F}$ is determined by T . That is, for every $\varphi \in \mathcal{F}$, either $\varphi \in T$ or $\neg\varphi \in T$.

Formulas φ and ψ are *(logically) equivalent*, written $\varphi \equiv \psi$, if $\varphi \vdash \psi$ and $\psi \vdash \varphi$. By (a), (iii), (iv), and the shorthand definition of \leftrightarrow , we find that $\varphi \equiv \psi$ if and only if $\varphi \leftrightarrow \psi \in \mathbf{Taut}$. Note that $\varphi \in \mathbf{Taut}$ if and only if $\varphi \equiv \top$. Also note that if $X \vdash \varphi$ and $\varphi \equiv \psi$, then $X \vdash \psi$.

More generally, if $X \subseteq \mathcal{F}$, we say that φ and ψ are *equivalent given X* , written $\varphi \equiv_X \psi$, if $X, \varphi \vdash \psi$ and $X, \psi \vdash \varphi$. As above, we have $\varphi \equiv_X \psi$ if and only if $\varphi \leftrightarrow \psi \in T(X)$. Also, $\varphi \in T(X)$ if and only if $\varphi \equiv_X \top$. Note that \equiv_\emptyset is simply \equiv . Also note that $\equiv_X \subseteq \equiv_{X'}$ whenever $X \subseteq X'$.

It can be shown that \equiv_X is a *congruence relation*, meaning it is an equivalence relation on \mathcal{F} such that

if $\varphi \equiv_X \varphi'$, then $\neg\varphi \equiv_X \neg\varphi'$, and

if C is countable and $\varphi_n \equiv_X \varphi'_n$ for $n \in C$, then $\bigwedge_{n \in C} \varphi_n \equiv_X \bigwedge_{n \in C} \varphi'_n$.

For $X, Y \subseteq \mathcal{F}$, we say that $X \equiv Y$ if $X \vdash Y$ and $Y \vdash X$. Note that $X \vdash Y$ if and only if $Y \subseteq T(X)$, which holds if and only if $T(Y) \subseteq T(X)$. Thus, $X \equiv Y$ if and only if $T(X) = T(Y)$. Also note that if $X \vdash \varphi$ and $X \equiv Y$, then $Y \vdash \varphi$.

The operations, \neg , \wedge , and \vee , pass in the usual way from \mathcal{F} to $B(X) = \mathcal{F}/\equiv_X$, making $B(X)$ into a Boolean σ -algebra, called the *Lindenbaum-Tarski σ -algebra of X* . If $[\varphi]_X \in B(X)$ denotes the equivalence class of φ , then $\varphi \equiv_X \psi$ if and only if $[\varphi]_X = [\psi]_X$. The partial order in $B(X)$ corresponds to the derivability relation. That is, $[\varphi]_X \leq [\psi]_X$ if and only if $X, \varphi \vdash \psi$. In $B(X)$, we have $0 = [\perp]_X$ and $1 = [\top]_X$.

We end this section with two items that we will need later. The first is a piece of notation. If T_0 and T_1 are theories with $T_0 \subseteq T_1$, then we write $[T_0, T_1]$ to denote the set of theories T that satisfy $T_0 \subseteq T \subseteq T_1$. The second is the following lemma.

(L:omit-psi) **Lemma 3.1.22.** *Let T be a theory, $\psi \in \mathcal{F}$, and $S \subseteq \mathcal{F}$. Define $S' = \{\psi \rightarrow \theta \mid \theta \in S\}$. Then $T + \psi + S = T + \psi + S'$.*

Proof. For each $\theta \in S$, we have $T + \psi + S \vdash \theta \vdash \psi \rightarrow \theta$. Thus, $T + \psi + S \vdash S'$, so that $T + \psi + S \vdash T + \psi + S'$. Conversely, for any $\theta \in S$, we have $\psi \rightarrow \theta, \psi \vdash \theta$, so that $S', \psi \vdash S$. Hence, $T + \psi + S' \vdash T + \psi + S$. \square

3.2 Inductive statements and entire sets

(S:entire) Let $\mathcal{F}^{\text{IS}} = \mathfrak{P}\mathcal{F} \times \mathcal{F} \times [0, 1]$. The elements, (X, φ, p) , of \mathcal{F}^{IS} are called *inductive statements*. Intuitively, we interpret (X, φ, p) as asserting that X partially entails φ , and that p is the degree of this partial entailment. In an inductive statement, X is called the *antecedent*, φ is called the *consequent*, and p is called the *probability*.

The remainder of this chapter is devoted to extending the derivability relation, \vdash , to inductive statements. Informally, the assertion, $Q \vdash (X, \varphi, p)$, where $Q \subseteq \mathcal{F}^{\text{IS}}$, means that (X, φ, p) can be derived from Q , using the rules of inductive inference. We will have nine such rules. They are:

- (R1) the rule of logical equivalence,
- (R2) the rule of logical implication,
- (R3) the rule of material implication,
- (R4) the rule of deductive transitivity,
- (R5) the addition rule,
- (R6) the multiplication rule,
- (R7) the continuity rule,
- (R8) the rule of inductive extension, and
- (R9) the rule of deductive extension.

The first rule, among other things, ensures that our collective inferences form a function mapping antecedent-consequent pairs, (X, φ) , to probabilities p . Rules (R2)–(R4) describe the relationship between deductive and inductive inference. Rules (R5)–(R7) are the usual mathematical rules for working with probabilistic assertions. And the final two rules provide a natural “completeness” to our inferences.

We follow the approach outlined in Remark 3.1.21. That is, we begin by defining an inductive theory, which will be a set of inductive statements that is closed under the nine rules of inductive inference, and satisfies certain connectivity requirements. This will allow us to speak of the inductive theory generated by a set $Q \subseteq \mathcal{F}^{\text{IS}}$, which we denote by $\mathbf{P}(Q)$. We then take $Q \vdash (X, \varphi, p)$ to mean that $(X, \varphi, p) \in \mathbf{P}(Q)$.

The notion of being closed under the nine rules of inductive inference will be built up in tiers. An admissible set is one that is closed under the first rule. An entire set is closed under the first seven rules. A semi-closed set is closed under the first eight rules. And a closed set is closed under all nine. In this section, we focus only on entire sets.

3.2.1 Seven of nine

We now formally state the first seven of the nine rules of inductive inference.

A set $P \subseteq \mathcal{F}^{\text{IS}}$ is *admissible* if it satisfies the *rule of logical equivalence*:

(R1) If $(X, \varphi, p) \in P$, $X' \equiv X$, and $\varphi' \equiv_X \varphi$, then $(X', \varphi', p) \in P$ and there is no other value p' such that $(X', \varphi', p') \in P$.

If P is admissible, then it is a function from $\mathfrak{P}\mathcal{F} \times \mathcal{F}$ to $[0, 1]$. In this case, we write $P(\varphi \mid X) = p$ to mean that $(X, \varphi, p) \in P$, and read the left-hand side, $P(\varphi \mid X)$, as *the probability of φ given X* . We also write X, ψ as shorthand for $X \cup \{\psi\}$, so that $P(\varphi \mid X, \psi)$ means $P(\varphi \mid X \cup \{\psi\})$. When $X = \emptyset$, we will omit it, leaving only $P(\varphi)$ or $P(\varphi \mid \psi)$. For admissible P , if $(X, \varphi, p) \in P$, then we say $P(\varphi \mid X)$ *exists* or is *defined*.

Note that any subset of an admissible set is also a function from $\mathfrak{P}\mathcal{F} \times \mathcal{F}$ to $[0, 1]$. We will therefore also use the notation $P(\varphi \mid X) = p$ for subsets of admissible sets.

If $P \subseteq \mathcal{F}^{\text{IS}}$, we define

$$\text{ante } P = \{X \subseteq \mathcal{F} : (X, \varphi, p) \in P \text{ for some } \varphi \in \mathcal{F} \text{ and } p \in [0, 1]\}.$$

That is, $X \in \text{ante } P$ if and only if X is the antecedent of some inductive statement in P .

The next six rules of inductive inference are encoded in the following definition.

Definition 3.2.1. A set $P \subseteq \mathcal{F}^{\text{IS}}$ is *entire* if it is admissible and satisfies the following:

(R2) (*the rule of logical implication*) If $X \in \text{ante } P$ and $X \vdash \varphi$, then $P(\varphi \mid X) = 1$.

(R3) (*the rule of material implication*) If $X \in \text{ante } P$ and $P(\psi \mid X, \varphi) = 1$, then $P(\varphi \rightarrow \psi \mid X) = 1$.

(R4) (*the rule of deductive transitivity*) If $P(\varphi \mid X) = 1$ and $\varphi \vdash \psi$, then $P(\psi \mid X) = 1$. Also, for any $X' \in \text{ante } P$, if $X' \vdash X$ and $P(\varphi \mid X) = 1$, then $P(\varphi \mid X') = 1$.

(R5) (*the addition rule*) Let $X \vdash \neg(\varphi \wedge \psi)$. Consider the equation,

$$P(\varphi \vee \psi \mid X) = P(\varphi \mid X) + P(\psi \mid X). \quad (3.2.1) \text{ add-rule}$$

If two of the above probabilities exist, then so does the third and (3.2.1) holds.

(R6) (*the multiplication rule*) Consider the equation,

$$P(\varphi \wedge \psi \mid X) = P(\varphi \mid X)P(\psi \mid X, \varphi). \quad (3.2.2) \text{ mult-rule}$$

If two of the above probabilities exist and are positive, then the third exists and (3.2.2) holds.

(R7) (*the continuity rule*) If $P(\varphi_n \mid X)$ exists and $X, \varphi_n \vdash \varphi_{n+1}$ for all $n \in \mathbb{N}$, then

$$P(\bigvee_n \varphi_n \mid X) = \lim_n P(\varphi_n \mid X). \quad (3.2.3) \text{ cont-rule}$$

Remark 3.2.2. If P is entire and $X \in \text{ante } P$, then X is consistent. To see this, suppose X is inconsistent. Choose $\varphi \in \mathcal{F}$ such that $X \vdash \varphi$ and $X \vdash \neg\varphi$. By the rule of logical implication $P(\varphi \mid X) = 1$ and $P(\neg\varphi \mid X) = 1$. Thus, by the addition rule, $P(\varphi \vee \neg\varphi) = 2$, which violates the definition of an inductive statement.

Remark 3.2.3. The first seven rules of inductive inference leave open the question of whether $P(\varphi \mid X) = 1$ implies $X \vdash \varphi$. In general, it does not, but a partial converse to the rule of logical implication will be given in Theorem 3.5.6.

3.2.2 Relative negation and certainty

(S:rel-neg-cert) Given an entire set P and a set $X \in \text{ante } P$, the domain of $P(\cdot \mid X)$ is not necessarily closed under conjunctions and disjunctions. It is, however, closed under relative negation. Also, conjunctions and disjunctions with a formula whose probability is 0 or 1 are still in the domain of $P(\cdot \mid X)$. These and related facts are described in this subsection.

(Expl:MathSEexpl) **Example 3.2.4.** The possible failure of the domain of $P(\cdot \mid X)$ to be closed under conjunctions and disjunctions can be seen in the following simple example, using $X = \emptyset$. Let $PV = \{\mathbf{r}_1, \mathbf{r}_2\}$. Fix $q \in (0, 1)$ and define $Q \subseteq \mathcal{F}^{\text{IS}}$ by $Q(\mathbf{r}_1) = Q(\mathbf{r}_2) = Q(\mathbf{r}_1 \leftrightarrow \mathbf{r}_2) = q$. In Example 4.3.8, we construct an entire set P such that $Q \subseteq P$, but in which $P(\mathbf{r}_1 \wedge \mathbf{r}_2)$ is undefined. As we will see in Theorem 3.2.18 below, this also implies $P(\mathbf{r}_1 \vee \mathbf{r}_2)$ is undefined.

$\langle \text{P:rel-neg} \rangle$ **Proposition 3.2.5.** *Let P be entire. If $P(\varphi \mid X)$ and $P(\psi \mid X)$ exist and $X, \varphi \vdash \psi$, then $P(\psi \wedge \neg\varphi \mid X)$ exists and*

$$P(\psi \wedge \neg\varphi \mid X) = P(\psi \mid X) - P(\varphi \mid X). \quad (3.2.4) \boxed{\text{rel-neg}}$$

In particular, $P(\varphi \mid X) \leq P(\psi \mid X)$.

Proof. Let $\psi' = \psi \wedge \neg\varphi$. Since $\psi \equiv_X \varphi \vee \psi'$, the rule of logical equivalence implies $P(\varphi \vee \psi' \mid X)$ exists. Since $\varphi \wedge \psi'$ is a contradiction, the addition rule implies that $P(\psi' \mid X)$ exists and

$$P(\varphi \vee \psi' \mid X) = P(\varphi \mid X) + P(\psi' \mid X),$$

which gives (3.2.4). \square

Remark 3.2.6. The final conclusion of Proposition 3.2.5 is referred to as the *monotonicity property* of P .

$\langle \text{C:rel-neg} \rangle$ **Corollary 3.2.7.** *Let P be entire. If $P(\varphi \mid X)$ exists, then $P(\neg\varphi \mid X)$ exists and*

$$P(\neg\varphi \mid X) = 1 - P(\varphi \mid X). \quad (3.2.5) \boxed{\text{prob-neg}}$$

Proof. Suppose $P(\varphi \mid X)$ exists. Then $X \in \text{ante } P$, so by the rule of logical implication, $P(\top \mid X) = 1$. Applying Proposition 3.2.5 with $\psi = \top$, and using $\top \wedge \neg\varphi \equiv \neg\varphi$ together with the rule of logical equivalence, we obtain (3.2.5). \square

$\langle \text{R:rel-neg} \rangle$ **Remark 3.2.8.** Proposition 3.2.5 requires neither the multiplication rule nor the continuity rule in its proof. Consequently, Corollary 3.2.7 also does not require them.

$\langle \text{P:cert-cl-conv} \rangle$ **Proposition 3.2.9.** *Let P be entire. If $P(\psi \mid X) = 1$ and $P(\varphi \wedge \psi \mid X)$ exists, then $P(\varphi \mid X) = P(\varphi \wedge \psi \mid X)$.*

Proof. Since $\psi \vdash \neg\varphi \vee \psi$, the rule of deductive transitivity gives $P(\neg\varphi \vee \psi \mid X) = 1$. By Corollary 3.2.7 and the rule of logical equivalence, $P(\varphi \wedge \neg\psi \mid X) = 0$. Hence, the result follows from the addition rule and the rule of logical equivalence. \square

$\langle \text{L:cond-exist} \rangle$ **Lemma 3.2.10.** *Let P be entire and suppose $P(\varphi \mid X)$ exists. Then $X \cup \{\varphi\} \in \text{ante } P$ if and only if $P(\varphi \mid X) > 0$.*

Proof. Suppose $P(\varphi \mid X) > 0$. Applying the multiplication rule with $\psi = \varphi$, and using $\varphi \wedge \varphi \equiv \varphi$ together with the rule of logical equivalence, we get $P(\varphi \mid X, \varphi) = 1$, which implies $X \cup \{\varphi\} \in \text{ante } P$.

Now suppose $P(\varphi \mid X) = 0$ and $X \cup \{\varphi\} \in \text{ante } P$. Then (3.2.5) implies $P(\neg\varphi \mid X) = 1$. Since $X, \varphi \vdash X$, the rule of deductive transitivity gives $P(\neg\varphi \mid X, \varphi) = 1$. Thus, again by (3.2.5), we have $P(\varphi \mid X, \varphi) = 0$. But this contradicts the rule of logical implication. \square

$\langle \text{P:certainty-closure} \rangle$ **Proposition 3.2.11.** *Let P be entire and suppose both $P(\varphi \mid X)$ and $P(\psi \mid X)$ exist. If $P(\varphi \mid X) \in \{0, 1\}$, then both $P(\varphi \vee \psi \mid X)$ and $P(\varphi \wedge \psi \mid X)$ exist, and*

$$\begin{aligned} P(\varphi \vee \psi \mid X) &= \max\{P(\varphi \mid X), P(\psi \mid X)\}, \\ P(\varphi \wedge \psi \mid X) &= \min\{P(\varphi \mid X), P(\psi \mid X)\}, \end{aligned}$$

Proof. By (3.2.5) and the rule of logical equivalence, it suffices to consider the case $\varphi \wedge \psi$. Suppose $P(\varphi \mid X) = 0$. By (3.2.5), we have $P(\neg\varphi \mid X) = 1$. Since $\neg\varphi \vdash \neg\varphi \vee \neg\psi$, the rule of deductive transitivity gives $P(\neg\varphi \vee \neg\psi \mid X) = 1$. By (3.2.5) and the rule of logical equivalence, $P(\varphi \wedge \psi \mid X) = 0$, proving the claim.

Now suppose $P(\varphi \mid X) = 1$. If $P(\psi \mid X) = 0$, then we are done by the previous case. Assume then that $P(\psi \mid X) > 0$. By Lemma 3.2.10, we have $X \cup \{\psi\} \in \text{ante } P$. Since $X, \psi \vdash X$, the rule of deductive transitivity gives $P(\varphi \mid X, \psi) = 1$. By the multiplication rule, $P(\psi \wedge \varphi \mid X) = P(\psi \mid X)$. \square

$\langle \text{C:certainty-closure} \rangle$ **Corollary 3.2.12.** *Let P be entire and suppose that $P(\varphi_n \mid X) = 1$ for all $n \in \mathbb{N}$. Then $P(\bigwedge_n \varphi_n \mid X) = 1$.*

Proof. By Proposition 3.2.11 and induction, we have $P(\bigwedge_1^n \varphi_j \mid X) = 1$ for all n . By (3.2.5) and the rule of logical equivalence, $P(\bigvee_1^n \neg\varphi_j \mid X) = 0$. The continuity rule then gives $P(\bigvee_n \neg\varphi_n \mid X) = 0$, which implies $P(\bigwedge_n \varphi_n \mid X) = 1$. \square

For $Q \subseteq \mathcal{F}^{\text{IS}}$ and $\mathcal{X} \subseteq \text{ante } Q$, let

$$\tau(Q; \mathcal{X}) = \{\theta \in \mathcal{F} \mid (X, \theta, 1) \in Q \text{ for all } X \in \mathcal{X}\}. \quad (3.2.6) \quad \boxed{\text{tau-Q-X}}$$

For $X \in \text{ante } Q$, we write $\tau(Q; X)$ for $\tau(Q; \{X\})$. We also write $\tau(Q)$, or τ_Q , for $\tau(Q; \text{ante } Q)$. Informally, $\tau(Q)$ is the set of all formulas that are true under Q , regardless of the antecedent used.

$\langle \text{P:tau-ded-theory} \rangle$ **Proposition 3.2.13.** *If P is entire and $\mathcal{X} \subseteq \text{ante } P$, then $\tau(P; \mathcal{X})$ is a deductive theory.*

Proof. Suppose $\tau(P; \mathcal{X}) \vdash \varphi$. By σ -compactness, choose countable $\Phi \subseteq \tau(P; \mathcal{X})$ such that $\Phi \vdash \varphi$. By Corollary 3.2.12, we have $\bigwedge \Phi \in \tau(P; \mathcal{X})$. Now let $X \in \mathcal{X}$. Then $P(\bigwedge \Phi \mid X) = 1$ and $\bigwedge \Phi \vdash \varphi$. By deductive transitivity, $P(\varphi \mid X) = 1$. Hence, $\varphi \in \tau(P; \mathcal{X})$, so $\tau(P; \mathcal{X})$ is a deductive theory. \square

3.2.3 Inductive vs. deductive inference

Rules (R1)–(R4), the rules of logical equivalence, logical implication, material implication, and deductive transitivity, describe the relationship between inductive and deductive inference. The next three results provide useful generalizations of these rules.

$\langle \text{P:log-equiv-gen} \rangle$ **Proposition 3.2.14.** *Let P be entire. If $P(\varphi \mid X) = p$ and $P(\varphi \leftrightarrow \varphi' \mid X) = 1$, then $P(\varphi' \mid X) = p$.*

Proof. Since $\varphi \wedge (\varphi \leftrightarrow \varphi') \equiv \varphi \wedge \varphi'$, the rule of logical equivalence together with Proposition 3.2.11 imply $P(\varphi \wedge \varphi' \mid X) = p$. Also, $\varphi \leftrightarrow \varphi' \vdash \varphi \vee \neg\varphi'$, so by the rule of deductive transitivity, $P(\varphi \vee \neg\varphi' \mid X) = 1$. Thus, (3.2.5) implies $P(\neg\varphi \wedge \varphi' \mid X) = 0$. Hence, by the addition rule and the rule of logical equivalence, $P(\varphi' \mid X) = p$. \square

$\langle P:\text{add-to-root} \rangle$ **Proposition 3.2.15.** *Let P be entire. If $P(\varphi \rightarrow \psi \mid X) = 1$ and $P(\varphi \mid X) > 0$, then $P(\psi \mid X, \varphi) = 1$.*

Proof. Suppose $P(\varphi \rightarrow \psi \mid X) = 1$ and $P(\varphi \mid X) > 0$. By (3.2.5) and the rule of logical equivalence, $P(\varphi \wedge \neg\psi \mid X) = 0$. Thus (3.2.4) implies $P(\varphi \wedge \psi \mid X) = P(\varphi \mid X) > 0$. Thus, by the multiplication rule, $P(\psi \mid X, \varphi) = 1$. \square

$\langle P:\text{transitivity} \rangle$ **Proposition 3.2.16.** *Let P be entire with $X \in \text{ante } P$ and $P(\varphi \mid X, \zeta) = 1$. Assume at least one of the following holds:*

- (i) $X, \varphi \vdash \psi$,
- (ii) $P(\psi \mid X, \varphi) = 1$,
- (iii) $P(\varphi \rightarrow \psi \mid X) = 1$.

Then $P(\psi \mid X, \zeta) = 1$.

Proof. Suppose $X \in \text{ante } P$ and $P(\varphi \mid X, \zeta) = 1$. First note that (i) is equivalent to $X \vdash \varphi \rightarrow \psi$. Thus, by the rule of logical implication, (i) implies (iii). Also, by the rule of material implication, (ii) implies (iii). Hence, we may assume (iii) holds. In this case, deductive transitivity gives $P(\varphi \rightarrow \psi \mid X, \zeta) = 1$. By (3.2.5) and the rule of logical equivalence, $P(\varphi \wedge \neg\psi \mid X, \zeta) = 0$. Thus, by (3.2.4), we have $P(\varphi \wedge \psi \mid X, \zeta) = 1$. Finally, since $\varphi \wedge \psi \vdash \psi$, deductive transitivity gives $P(\psi \mid X, \zeta) = 1$. \square

3.2.4 Generalizations of the addition rule

$\langle L:\text{fin-add} \rangle$ **Lemma 3.2.17.** *Let P be entire and suppose $X \vdash \neg(\varphi_i \wedge \varphi_j)$ whenever $i \neq j$. If $P(\varphi_i \mid X)$ exists for $1 \leq i \leq n$, then $P(\bigvee_{i=1}^n \varphi_i) = \sum_{i=1}^n P(\varphi_i \mid X)$.*

Proof. The case $n = 2$ is the addition rule. Suppose the lemma holds for some $n = k$ and consider the case $n = k + 1$. Note that $X \vdash \neg(\varphi \wedge \psi)$ is equivalent to $X, \psi \vdash \neg\varphi$. Thus, $X, \varphi_{k+1} \vdash \neg\varphi_i$ for all $i \leq k$. This implies $X, \varphi_{k+1} \vdash \bigwedge_1^k \neg\varphi_i \equiv \neg\bigvee_1^k \varphi_i$. Therefore, $X \vdash \neg(\bigvee_1^k \varphi_i \wedge \varphi_{k+1})$, so the addition rule and the inductive hypothesis show that the result holds for $n = k + 1$. \square

$\langle T:\text{incl-excl} \rangle$ **Theorem 3.2.18 (Inclusion-exclusion).** *Let P be entire and consider the equation*

$$P(\varphi \vee \psi \mid X) = P(\varphi \mid X) + P(\psi \mid X) - P(\varphi \wedge \psi \mid X). \quad (3.2.7) \quad \boxed{\text{incl-excl}}$$

If three of the above probabilities exist, then so does the fourth and (3.2.7) holds.

Proof. Let $\zeta_1 = \varphi \wedge \neg\psi$ and $\zeta_2 = \neg\varphi \wedge \psi$. Suppose three of the probabilities in (3.2.7) exist. There are four possible cases.

The first case is that $P(\varphi | X)$, $P(\psi | X)$, and $P(\varphi \wedge \psi | X)$ exist. In this case, since $\zeta_1 \equiv \varphi \wedge \neg(\varphi \wedge \psi)$, Proposition 3.2.5 implies that $P(\zeta_1 | X)$ exists. Similarly, $P(\zeta_2 | X)$ exists. Thus, by the addition rule and the fact that $\varphi \vee \psi \equiv \zeta_1 \vee \zeta_2 \vee (\varphi \wedge \psi)$, we have that $P(\varphi \vee \psi | X)$ exists.

The second case is that $P(\varphi \vee \psi | X)$, $P(\varphi | X)$, and $P(\psi | X)$ exist. In this case, since $\zeta_1 \equiv (\varphi \vee \psi) \wedge \neg\psi$, Proposition 3.2.5 implies that $P(\zeta_1 | X)$ exists. Similarly, $P(\zeta_2 | X)$ exists. By the addition rule, Proposition 3.2.5, and the fact that $\varphi \wedge \psi \equiv (\varphi \vee \psi) \wedge \neg(\zeta_1 \vee \zeta_2)$, we have that $P(\varphi \wedge \psi | X)$ exists.

By symmetry, the third and fourth cases are covered by the assumption that $P(\varphi \vee \psi | X)$, $P(\varphi | X)$, and $P(\varphi \wedge \psi | X)$ exist. The argument from the second case shows that $P(\zeta_2 | X)$ exists. By the addition rule and the fact that $\psi \equiv (\varphi \wedge \psi) \vee \zeta_2$, we have that $P(\psi | X)$ exists.

Hence, in all cases, all four probabilities in (3.2.7) exist. By Lemma 3.2.17,

$$\begin{aligned} P(\varphi | X) &= P(\zeta_1 | X) + P(\varphi \wedge \psi | X), \\ P(\psi | X) &= P(\zeta_2 | X) + P(\varphi \wedge \psi | X), \text{ and} \\ P(\varphi \vee \psi | X) &= P(\zeta_1 | X) + P(\zeta_2 | X) + P(\varphi \wedge \psi | X). \end{aligned}$$

Putting these together yields (3.2.7). \square

Remark 3.2.19. By Theorem 3.2.18, if P is entire, then in the addition rule, it is not necessary that $X \vdash \neg(\varphi \wedge \psi)$. It is sufficient that $P(\varphi \wedge \psi | X) = 0$.

3.2.5 Generalizations of the multiplication rule

Theorem 3.2.20. *If P is entire, then in the multiplication rule, it is not necessary that the two defined probabilities be positive. It is enough to assume that solving for the third probability does not result in dividing by zero.*

Proof. First suppose $P(\varphi | X)$ and $P(\psi | X, \varphi)$ both exist. By Lemma 3.2.10, we have $P(\varphi | X) > 0$. Suppose $P(\psi | X, \varphi) = 0$. Then $P(\neg\psi | X, \varphi) = 1$, by (3.2.5), so by the multiplication rule, $P(\varphi \wedge \neg\psi | X) = P(\varphi | X)$. Therefore, Proposition 3.2.5 implies $P(\varphi \wedge \psi | X) = 0$.

Next, suppose $P(\varphi | X) > 0$ and $P(\varphi \wedge \psi | X) = 0$. Then Proposition 3.2.5 implies $P(\varphi \wedge \neg\psi | X) = P(\varphi | X) > 0$, so by the multiplication rule, $P(\neg\psi | X, \varphi) = 1$. Applying Proposition 3.2.5 again gives $P(\psi | X, \varphi) = 0$.

Finally, suppose $P(\psi | X, \varphi) > 0$ and $P(\varphi \wedge \psi | X)$ exists. By Lemma 3.2.10, we have $X \cup \{\varphi, \psi\} \in \text{ante } P$. But $X \cup \{\varphi, \psi\} \equiv X \cup \{\varphi \wedge \psi\}$, so by the rule of logical equivalence, $X \cup \{\varphi \wedge \psi\} \in \text{ante } P$. Thus, by Lemma 3.2.10, we have $P(\varphi \wedge \psi | X) > 0$. \square

Theorem 3.2.21 (Bayes' theorem). *If P is entire, then*

$$P(\varphi | X)P(\psi | X, \varphi) = P(\psi | X)P(\varphi | X, \psi), \quad (3.2.8) \quad \boxed{\text{Bayes}}$$

provided that all four of the above probabilities exist.

Proof. Since $\varphi \wedge \psi \equiv \psi \wedge \varphi$, this follows immediately from the multiplication rule. \square

3.2.6 Generalizations of the continuity rule

Proposition 3.2.22. *Let P be entire. If $P(\varphi_n | X)$ exists and $X, \varphi_{n+1} \vdash \varphi_n$ for all $n \in \mathbb{N}$, then*

$$P(\bigwedge_n \varphi_n | X) = \lim_n P(\varphi_n | X).$$

Proof. Let $\psi_n = \neg\varphi_n$. By the continuity rule,

$$P(\bigvee_n \psi_n | X) = \lim_n P(\psi_n | X).$$

But $\bigvee_n \psi_n \equiv \neg \bigwedge_n \varphi_n$, so (3.2.5) gives the desired result. \square

(T:cont-rule) **Theorem 3.2.23.** *Let P be entire and assume $P(\varphi_n | X)$ exists for all n . If $P(\varphi_{n+1} | X, \varphi_n) = 1$ for all n , then*

$$P(\bigvee_n \varphi_n | X) = \lim_n P(\varphi_n | X). \quad (3.2.9) \text{ cont-from-below}$$

Similarly, if $P(\varphi_n | X, \varphi_{n+1}) = 1$ for all n , then

$$P(\bigwedge_n \varphi_n | X) = \lim_n P(\varphi_n | X). \quad (3.2.10) \text{ cont-from-above}$$

Proof. Suppose $P(\varphi_{n+1} | X, \varphi_n) = 1$ for all n . By Lemma 3.2.10, we have $P(\varphi_n | X) > 0$ for all n . Let $\psi_n = \bigvee_{j=1}^n \varphi_j$. We first show that

$$P(\psi_n | X) = P(\varphi_n | X) > 0, \text{ and} \quad (3.2.11) \text{ cont-1}$$

$$P(\varphi_{n+1} | X, \psi_n) = 1, \quad (3.2.12) \text{ cont-2}$$

for all n .

Since $\psi_1 = \varphi_1$, we have that (3.2.11) and (3.2.12) hold for $n = 1$. Suppose it is true for $n - 1$. Then

$$\psi_n \equiv \psi_{n-1} \vee \varphi_n \equiv \varphi_n \vee (\psi_{n-1} \wedge \neg\varphi_n).$$

But $P(\varphi_n | X, \psi_{n-1}) = 1$, so it follows from the rule of material implication and (3.2.5) that $P(\psi_{n-1} \wedge \neg\varphi_n | X) = 0$. Hence, Proposition 3.2.11 gives us (3.2.11).

By Proposition 3.2.16 and induction, we have $P(\varphi_{n+1} | X, \varphi_m) = 1$ for all $m \leq n$, which gives $P(\varphi_m \rightarrow \varphi_{n+1} | X) = 1$ for all $m \leq n$ by the rule of material implication. Note that $\bigwedge_{m=1}^n (\varphi_m \rightarrow \varphi_{n+1}) \equiv \psi_n \rightarrow \varphi_{n+1}$. Hence, by Proposition 3.2.11 and induction, we have $P(\psi_n \rightarrow \varphi_{n+1} | X) = 1$. Since $P(\psi_n | X) > 0$, Proposition 3.2.15 gives (3.2.12).

Having established (3.2.11) and (3.2.12) for all n , observe that $\psi_n \vdash \psi_{n+1}$ for all n . Thus, by the continuity rule,

$$P(\bigvee_n \psi_n | X) = \lim_n P(\psi_n | X).$$

Using (3.2.11) and the fact that $\bigvee_n \psi_n \equiv \bigvee_n \varphi_n$, we obtain (3.2.9).

For (3.2.10), assume $P(\varphi_n \mid X, \varphi_{n+1}) = 1$ for all n . Define $\psi_n = \neg\varphi_n$. By the rule of material implication, $P(\varphi_{n+1} \rightarrow \varphi_n \mid X) = 1$ for all n , which gives $P(\psi_n \rightarrow \psi_{n+1} \mid X) = 1$ for all n .

We first suppose that $P(\varphi_n \mid X) < 1$ for all n . Since $P(\psi_n \mid X) > 0$ by (3.2.5), Proposition 3.2.15 implies $P(\psi_{n+1} \mid X, \psi_n) = 1$ for all n . Applying (3.2.9) and (3.2.5) gives (3.2.10).

Now suppose that $P(\varphi_n \mid X) = 1$ for some n . By Proposition 3.2.16, we have $P(\varphi_j \mid X) = 1$ for all $j \leq n$. Let $n_0 = \sup\{n : P(\varphi_n \mid X) = 1\}$. If $n_0 = \infty$, then $P(\varphi_n \mid X) = 1$ for all n . This implies $P(\psi_n \mid X) = 0$ for all n by Proposition 3.2.11. Hence, by Corollary 3.2.12 and (3.2.5), we have $P(\bigvee_n \psi_n \mid X) = 0$, and this establishes (3.2.10). Assume, then, that $n_0 < \infty$, so that $P(\varphi_n \mid X) = 1$ for all $n \leq n_0$ and $P(\varphi_n \mid X) < 1$ for all $n > n_0$. By what we have already proven,

$$P(\bigwedge_{n_0+1}^{\infty} \varphi_n \mid X) = \lim_n P(\varphi_n \mid X).$$

By Proposition 3.2.11 and induction, we have $P(\bigwedge_n \varphi_n \mid X) = P(\bigwedge_{n_0+1}^{\infty} \varphi_n \mid X)$, and so (3.2.10) holds. \square

(T:ctbl-add) **Theorem 3.2.24 (Countable additivity).** *If P is entire and $P(\varphi_i \wedge \varphi_j \mid X) = 0$ for all $1 \leq i < j < \infty$, then*

$$P(\bigvee_n \varphi_n \mid X) = \sum_n P(\varphi_n \mid X).$$

Proof. Let $\psi_n = \bigvee_1^n \varphi_j$, so that $\psi_n \vdash \psi_{n+1}$ and $\bigvee_n \psi_n \equiv \bigvee_n \varphi_n$. Note that $\psi_n \wedge \varphi_{n+1} \equiv \bigvee_1^n (\varphi_j \wedge \varphi_{n+1})$. By Proposition 3.2.11 and induction, we have $P(\psi_n \wedge \varphi_{n+1} \mid X) = 0$. Thus, by Theorem 3.2.18,

$$P(\psi_n \vee \varphi_{n+1} \mid X) = P(\psi_n \mid X) + P(\varphi_{n+1} \mid X).$$

It follows by induction that $P(\psi_n \mid X) = \sum_1^n P(\varphi_j \mid X)$ for all n . Letting $n \rightarrow \infty$ and applying the continuity rule completes the proof. \square

3.3 Closed sets and inductive derivability

(S:closed) 3.3.1 The rule of inductive extension

The first seven rules of inductive inference encapsulate the core of our inductive calculus. There are, however, two important and essential supplemental rules we must define. The first is called the “rule of inductive extension.” To motivate this rule, recall the situation in Example 3.2.4. As mentioned therein, we will later construct an entire set P in which $P(\mathbf{r}_1) = P(\mathbf{r}_2) = P(\mathbf{r}_1 \leftrightarrow \mathbf{r}_2) = q$, but $P(\mathbf{r}_1 \wedge \mathbf{r}_2)$ is undefined (see Example 4.3.8). This situation is entirely satisfactory and will not violate our rules of inference in any way. It is, in fact, self-evident that without additional information, there is no way to deduce a probability for $\mathbf{r}_1 \wedge \mathbf{r}_2$ based solely on the probabilities of \mathbf{r}_1 and \mathbf{r}_2 .

We can, however, decide to assign, a priori, a probability to $\mathbf{r}_1 \wedge \mathbf{r}_2$. In other words, given a value q' , we may wish to consider the set $P' = P \cup \{(\emptyset, \mathbf{r}_1 \wedge \mathbf{r}_2, q')\}$. Of course, we must choose q' so that our new, enlarged set P' continues to conform to the seven rules of inference we have already established. The question naturally arises: which values of q' are possible?

In a situation such as this, one of three things can occur. The first is illustrated by the case $q = 1/4$. In this case, we show in Example 4.3.10 that there are no possible values of q' . In other words, although P is entire, there is no way to assign a probability to $\mathbf{r}_1 \wedge \mathbf{r}_2$ without violating one of our first seven rules. There is, therefore, something defective about P , but this flaw cannot be seen from our first seven rules alone.

The second possibility is illustrated by the case $q = 1/2$. In this case, we show in Proposition 4.3.11 that $q' = 1/4$ is the unique value that works. In other words, the only way to assign a probability to $\mathbf{r}_1 \wedge \mathbf{r}_2$ without violating one of our first seven rules is to assign it probability $1/4$. In this case, it seems reasonable that the uniqueness of this value ought to let us infer $(\emptyset, \mathbf{r}_1 \wedge \mathbf{r}_2, 1/4)$ from P . But such an inference is not possible with only the first seven rules, because P is already entire.

The third possibility is that there are multiple values of q' that work. Although this case does not arise in Example 3.2.4, it is clear that it can arise in even simpler examples. In a case such as this, there is nothing necessarily defective about our entire set P , but at the same time, we cannot make any inference about the probability of $\mathbf{r}_1 \wedge \mathbf{r}_2$.

To describe the rule that will rectify these situations, we begin by defining a new kind of set, which we call “complete.” Even after our inductive calculus is fully developed, the process of inductive inference will not typically produce complete sets. Rather, they represent a sort of ideal in which all meaningfully connected probabilities have been logically determined.

(D:complete) **Definition 3.3.1.** A set $\bar{P} \subseteq \mathcal{F}^{\text{IS}}$ is *complete* if it is entire and satisfies the following conditions:

- (i) If $\bar{P}(\varphi | X)$ and $\bar{P}(\psi | X)$ exist, then $\bar{P}(\varphi \wedge \psi | X)$ exists.
- (ii) If $X \in \text{ante } \bar{P}$ and $X \cup \{\varphi\} \in \text{ante } \bar{P}$, then $\bar{P}(\varphi | X)$ exists.

Remark 3.3.2. In general, entire sets obey neither (i) nor (ii) in the definition above. Example 3.2.4 describes an entire set that violates (i). In Example 4.3.12, we exhibit an entire set P with $P(\mathbf{r}_1) = 1/2$ and $P(\mathbf{r}_2 | \mathbf{r}_3) = 1$, but with $P(\mathbf{r}_3)$ undefined, thereby violating (ii) with $X = \emptyset$ and $\varphi = \mathbf{r}_3$.

Having defined complete sets, we are now in a position to state our eighth rule of inductive inference. A set which is closed under the first eight rules will be called “semi-closed.”

If $P \subseteq \bar{P} \subseteq \mathcal{F}^{\text{IS}}$, then \bar{P} is called an *extension* of P . A complete extension will also be called a *completion*. A set $P \subseteq \mathcal{F}^{\text{IS}}$ is *semi-closed* if it is entire and satisfies the *rule of inductive extension*:

(R8) If $\overline{P}(\varphi \mid X) = p$ for every completion \overline{P} of P , then $P(\varphi \mid X) = p$.

Note that every complete set is semi-closed.

Every semi-closed set has a completion. To see this, suppose P is semi-closed but has no completion. Then (R8) implies $P = \mathcal{F}^{\text{IS}}$. But then P is not admissible, and therefore not entire, a contradiction.

This means that an entire set can fail to be semi-closed in two ways. On the one hand, it can simply not have enough probabilities. This happens, for instance, in Example 3.2.4 when $q = 1/2$. In this case, our set P will not be semi-closed until we add more probabilities, including $(\emptyset, \mathbf{r}_1 \wedge \mathbf{r}_2, 1/4)$, as required by (R8).

On the other hand, it can fail to be semi-closed because it has no completion. This happens, for instance, in Example 3.2.4 when $q = 1/4$, since in that case there is no way to assign a probability to $\mathbf{r}_1 \wedge \mathbf{r}_2$ without breaking the entirety of the set.

The following result is an analogue of Proposition 3.2.14 for antecedents, but it requires that P be semi-closed.

(P:log-equiv-gen-2) **Proposition 3.3.3.** *Let P be semi-closed. If $P(\varphi \mid X, \psi) = p$ and $P(\psi \leftrightarrow \psi' \mid X) = 1$, then $P(\varphi \mid X, \psi') = p$.*

Proof. Let P be semi-closed with $P(\varphi \mid X, \psi) = p$ and $P(\psi \leftrightarrow \psi' \mid X) = 1$. Let \overline{P} be a completion of P . Then $\overline{P}(\varphi \mid X, \psi) = p$ and $\overline{P}(\psi \leftrightarrow \psi' \mid X) = 1$. By Definition 3.3.1(ii), we have that $\overline{P}(\psi \mid X) = q$ for some $q \in [0, 1]$, and Lemma 3.2.10 implies $q > 0$. By the multiplication rule, $\overline{P}(\varphi \wedge \psi \mid X) = pq$.

By Proposition 3.2.11, we have $\overline{P}(\varphi \wedge \psi \wedge (\psi \leftrightarrow \psi') \mid X) = pq$. By the rule of logical equivalence, $\overline{P}(\varphi \wedge \psi' \wedge (\psi \leftrightarrow \psi') \mid X) = pq$. Proposition 3.2.9 then implies $\overline{P}(\varphi \wedge \psi' \mid X) = pq$. By Proposition 3.2.14, we have $\overline{P}(\psi' \mid X) = q$. Hence, by the multiplication rule, $\overline{P}(\varphi \mid X, \psi') = p$. Since \overline{P} was arbitrary, the rule of inductive extension gives $P(\varphi \mid X, \psi') = p$. \square

3.3.2 The rule of deductive extension

Our final rule is the “rule of deductive extension.” Informally, it says that any antecedent can be freely expanded to include any number of formulas already known to have probability one. A set of inductive statements that is closed under all nine rules of inference will be called “closed.” More specifically, a semi-closed set P is said to be *closed* if it satisfies the *rule of deductive extension*:

(R9) If $S \subseteq \mathcal{F}$ is nonempty and $P(\theta \mid X) = 1$ for all $\theta \in S$, then $X \cup S \in \text{ante } P$ and $P(\cdot \mid X, S) = P(\cdot \mid X)$.

With this final definition, our rules of inductive inference are complete.

3.3.3 Pre-theories

We now wish to use these rules to define inductive derivability. Our aim is to make sense of the statement $Q \vdash (X, \varphi, p)$. Informally, we imagine $Q \vdash (X, \varphi, p)$

to mean that, starting from the inductive statements in Q , we may use a sequence of applications of our nine rules of inference to derive the inductive statement, (X, φ, p) . In keeping with the spirit of our infinitary language \mathcal{F} , we will imagine, when necessary, that this sequence is at most countable. We aim to make this notion precise by using our nine rules of inference and their related closure properties (admissible, entire, semi-closed, and closed).

As mentioned earlier, we plan to follow the route described in Remark 3.1.21. That is, by analogy with deductive theories, we want to define an “inductive theory.” We will then say that $Q \vdash (X, \varphi, p)$ if (X, φ, p) is an element of the smallest inductive theory containing Q . Our first task, of course, is to determine exactly what ought to constitute an inductive theory. At first glance, the answer may seem trivially obvious: an inductive theory ought to simply be a closed set. After all, closed sets, by definition, are those sets that are closed under all nine rules of inference, so this would be the natural analogue of a deductive theory. We will see, however, that closed subsets of \mathcal{F}^{IS} can be much larger than we might initially expect, and as such, do not fit our intuitive understanding of what an inductive theory ought to be.

To see this, let us first focus our attention on semi-closed sets, which are closed under (R1)–(R8). Imagine we have a starting collection of inductive statements, Q . Using Q together with rules (R1)–(R8), we begin making inferences and adding new inductive statements to our collection. When we have exhausted all inferences that are possible with (R1)–(R8), we arrive at a finalized collection, P_0 , which we will call a “pre-theory.” Our pre-theory P_0 is closed under (R1)–(R8), so by definition, P_0 is a semi-closed set. But P_0 has another important property that arbitrary semi-closed sets do not share. The elements of P_0 are all “connected” to the elements of Q in a certain sense. They are connected via (R1)–(R8).

To clarify the nature of this connection, let us consider the effects of (R1)–(R8) on the antecedents of Q . That is, imagine we use a single application of one of the rules (R1)–(R8) to infer (X, φ, p) from Q . We wish to know how X is related to ante Q . If we have used any of the rules (R2), (R3), (R4), (R5), (R7), or (R8), then we know that $X \in \text{ante } Q$. If we used (R1), then $X \equiv Y$ for some $Y \in \text{ante } Q$. And if we used (R6), then either $X \in \text{ante } Q$ or $X = Y \cup \{\varphi\}$ for some $Y \in \text{ante } Q$ and some $\varphi \in \mathcal{F}$. Generally speaking, no matter which of (R1)–(R8) we used, we may conclude that either $X \equiv Y$ or $X \equiv Y \cup \{\varphi\}$ for some $Y \in \text{ante } Q$ and $\varphi \in \mathcal{F}$. More generally, if (X, φ, p) is inferred from Q via a countable sequence of applications of (R1)–(R8), then $X \equiv Y \cup \Phi$ for some $Y \in \text{ante } Q$ and some countable $\Phi \subseteq \mathcal{F}$.

Motivated by this, we make the following definition. A set $Q \subseteq \mathcal{F}^{\text{IS}}$ is *strongly connected* if there exists $X_0 \in \text{ante } Q$ such that every $X \in \text{ante } Q$ is *countably axiomatizable* over X_0 . That is, for every $X \in \text{ante } Q$, there exists a countable set $\Phi \subseteq \mathcal{F}$ such that $X \equiv X_0 \cup \Phi$. Note that a strongly connected set is necessarily nonempty. A set $P_0 \subseteq \mathcal{F}^{\text{IS}}$ is a *pre-theory* if it is semi-closed and strongly connected.

Strong connectivity formalizes the notion that the inductive statements in a set can be related to one another via the calculus of (R1)–(R8). A pre-theory

represents the results of exhausting all possible inferences using (R1)–(R8). We require not only that a pre-theory be closed under (R1)–(R8), and therefore a semi-closed set, but also that it be strongly connected. A set which is semi-closed but not strongly connected will not violate (R1)–(R8), but it will include unnecessary parts and statements which can never be related to one another by the calculus of (R1)–(R8).

3.3.4 Inductive theories

(S:ind-theories) We now turn our attention to (R9). The first issue to address is the interplay between (R9) and the previous eight rules. Suppose we have exhausted all inferences from (R1)–(R8) and obtained a pre-theory, P_0 . We then use (R9) to infer a new inductive statement, (X, φ, p) . Is it possible that $P_0 \cup \{(X, \varphi, p)\}$ can be extended to an even larger pre-theory? Corollary 3.3.5 assures us that it cannot. In other words, we lose nothing by requiring that all applications of (R9) take place after all applications of (R1)–(R8).

With this in mind, we aim to say that an “inductive theory” is what we obtain by first constructing a pre-theory and then “closing it up” with (R9). Theorem 3.3.4 below shows that this closure operation is well-defined and produces a unique result. The proof of Theorem 3.3.4, as well as all the other results in the remainder of this section, will be postponed until Sections 3.4 and 3.5. We present here only the statements of the results, so that we may first see an overview of the entire construction.

(T:theory-defn) **Theorem 3.3.4.** *Every pre-theory has a unique smallest closed extension.*

If $P_0 \subseteq \mathcal{F}^{\text{IS}}$ is a pre-theory, let $\mathbf{P}(P_0)$ or \mathbf{P}_{P_0} denote its smallest closed extension.

(C:theory-defn) **Corollary 3.3.5.** *Let $P_0, P'_0 \subseteq \mathcal{F}^{\text{IS}}$ be pre-theories. If $\mathbf{P}(P_0) = \mathbf{P}(P'_0)$, then $P_0 = P'_0$. In particular, if P_0 is a pre-theory, then there is no pre-theory P'_0 such that $P_0 \subset P'_0 \subseteq \mathbf{P}(P_0)$.*

Having established these results, we will be able to say that an “inductive theory” is a set of the form $\mathbf{P}(P_0)$ for some pre-theory P_0 . Before formally defining it as such, we pause to characterize such sets in a way analogous to our characterization of pre-theories. For this characterization, we first define a new connectivity property.

Recall the notation established in (3.2.6). We say that a set $Q \subseteq \mathcal{F}^{\text{IS}}$ is *connected* if there exists a strongly connected $\widehat{Q} \subseteq Q$ such that for all $X \in \text{ante } Q$, there exists an $\widehat{X} \in \text{ante } \widehat{Q}$ and a set $S \subseteq \tau(\widehat{Q}; \widehat{X})$ such that $X \equiv \widehat{X} \cup S$. Any such \widehat{Q} will be called a *basis* for Q . In other words, a connected set is a “lift” of a strongly connected set, where we lift up the antecedents by including formulas that have probability one.

Note that if Q is strongly connected, then Q is connected and is its own basis. Also note the following important difference between connectivity and strong connectivity. Strong connectivity is a property of $\text{ante } Q$. That is, if

$\text{ante } Q = \text{ante } Q'$, then Q is strongly connected if and only if Q' is strongly connected. Connectivity, on the other hand, is not. Connectivity depends not only on $\text{ante } Q$, but also on $\{(X, \theta) \mid (X, \theta, 1) \in Q\}$.

$\langle \text{T:theory-char} \rangle$ **Theorem 3.3.6.** *Let $P \subseteq \mathcal{F}^{\text{IS}}$. The following are equivalent:*

- (i) $P = \mathbf{P}(P_0)$ for some (unique) pre-theory P_0 .
- (ii) P is closed and connected.

With this last result, we can finally state the definition that is the linchpin of our entire inductive calculus: a set $P \subseteq \mathcal{F}^{\text{IS}}$ is a *inductive theory* if it is closed and connected.

The intuitive interpretation of connectivity is analogous to strong connectivity, but instead of using only (R1)–(R8), we use the whole of our inductive calculus. That is, a connected set P is one whose inductive statements can be related to one another via the calculus. A set which is closed but not connected will not violate the calculus, but it will have unnecessary parts which the inductive calculus can never reach.

3.3.5 Inductive derivability

$\langle \text{S:ind-derivability} \rangle$ It is now straightforward to define inductive derivability. We begin by defining a set $Q \subseteq \mathcal{F}^{\text{IS}}$ to be *consistent* if it is connected and can be extended to an inductive theory. The requirement that a consistent set be extendable to an inductive theory ensures that it does not violate the calculus of inductive inference. The requirement that it be connected ensures that its statements are all logically related to one another.

Note that every pre-theory is consistent. Moreover, if P_0 is a pre-theory, then $\mathbf{P}(P_0)$ is the smallest extension of P_0 to an inductive theory.

$\langle \text{T:theory-gen-defn} \rangle$ **Theorem 3.3.7.** *Every consistent set has a unique smallest extension to an inductive theory.*

If $Q \subseteq \mathcal{F}^{\text{IS}}$ is consistent, let $\mathbf{P}(Q)$ or \mathbf{P}_Q denote its smallest extension to an inductive theory. We call \mathbf{P}_Q the *inductive theory generated by Q* . If $Q \subseteq \mathcal{F}^{\text{IS}}$ and $(X, \varphi, p) \in \mathcal{F}^{\text{IS}}$, we write $Q \vdash (X, \varphi, p)$ to mean that Q is consistent and $\mathbf{P}_Q(\varphi \mid X) = p$. When the turnstile symbol, \vdash , is used in this fashion, we will call it the *inductive derivability relation*. When $Q \vdash (X, \varphi, p)$, we say that (X, φ, p) is *inductively derivable* from Q , or that Q *proves* (X, φ, p) . Note that unlike deductive derivability, our convention is that if $Q \subseteq \mathcal{F}^{\text{IS}}$ is inconsistent, then Q does not prove anything.

$\langle \text{R:classic-ind-th-defn} \rangle$ **Remark 3.3.8.** If $P \subseteq \mathcal{F}^{\text{IS}}$ is consistent, then P is an inductive theory if and only if $P = \mathbf{P}(P)$, which holds if and only if $\mathbf{P}(P) \subseteq P$. Hence, using the above definition of inductive derivability, we can say that a consistent set $P \subseteq \mathcal{F}^{\text{IS}}$ is an inductive theory if and only if $P \vdash (X, \varphi, p)$ implies $(X, \varphi, p) \in P$ for all $(X, \varphi, p) \in \mathcal{F}^{\text{IS}}$.

3.4 Connectivity, restrictions, and lifts

(S:lifts) In this section, we prove the results in Section 3.3.4. We begin with four subsections on preliminary results needed in the proofs. In the first two subsections, we establish some basic facts about connected and strongly connected sets, and how they relate to the rules of inference. The third subsection looks at restrictions of sets $Q \subseteq \mathcal{F}^{\text{IS}}$, and examines which closure properties are preserved under restriction. Finally, the fourth subsection defines the “lift” of a pre-theory, which we denote by $\mathbf{L}(P_0)$. After establishing these preliminaries, we give the proofs of Theorem 3.3.4, Corollary 3.3.5, and Theorem 3.3.6.

3.4.1 Connectivity properties

For $X, X_0 \subseteq \mathcal{F}^{\text{IS}}$, we write $X \hookrightarrow X_0$ to mean that X is countably axiomatizable over X_0 . If $\mathcal{X} \subseteq \mathfrak{P}\mathcal{F}$, we write $X \hookrightarrow \mathcal{X}$ to mean $X \hookrightarrow X_0$ for some $X_0 \in \mathcal{X}$.

(P:ctbly-ax=one-ax) **Proposition 3.4.1.** *Let $X, X_0 \subseteq \mathcal{F}$. Then $X \hookrightarrow X_0$ if and only if $X \equiv X_0 \cup \{\psi\}$ for some $\psi \in \mathcal{F}$.*

Proof. This follows immediately from the fact that $\Phi \equiv \bigwedge \Phi$ for any countable $\Phi \subseteq \mathcal{F}$. \square

(P:root-exist) **Proposition 3.4.2.** *If $Q \subseteq \mathcal{F}^{\text{IS}}$ is connected, then there exists $X_0 \in \text{ante } Q$ such that $X \vdash X_0$ for all $X \in \text{ante } Q$. This X_0 is unique in the sense that if X'_0 is another such antecedent, then $X_0 \equiv X'_0$.*

Proof. Let \widehat{Q} be a basis for Q . Since \widehat{Q} is strongly connected, we may choose $X_0 \in \text{ante } \widehat{Q}$ such that $\widehat{X} \hookrightarrow X_0$ for all $\widehat{X} \in \text{ante } \widehat{Q}$. Now let $X \in \text{ante } Q$ be given. Choose $\widehat{X} \in \text{ante } \widehat{Q}$ and $S \subseteq \tau(\widehat{Q}; \widehat{X})$ such that $X \equiv \widehat{X} \cup S$. Since $\widehat{X} \hookrightarrow X_0$, we have $\widehat{X} \vdash X_0$, and therefore $X \vdash X_0$. Uniqueness is immediate since $X_0 \vdash X'_0$ and $X'_0 \vdash X_0$ implies $X_0 \equiv X'_0$. \square

Let $Q \subseteq \mathcal{F}^{\text{IS}}$ be connected. Choose X_0 as in Proposition 3.4.2 and let $T_0 = T(X_0)$. By Proposition 3.4.2, the deductive theory T_0 does not depend on the choice of X_0 . We call T_0 the *root* of Q . Note that if Q is admissible, then by the rule of logical equivalence, $T_0 \in \text{ante } Q$.

(P:str-conn-root-one-away) **Proposition 3.4.3.** *If $Q \subseteq \mathcal{F}^{\text{IS}}$ is strongly connected with root T_0 , then for each $X \in \text{ante } Q$, there exists $\psi \in \mathcal{F}$ such that $T(X) = T_0 + \psi$.*

Proof. Since Q is strongly connected, we may choose $X_0 \in \text{ante } Q$ such that $X \hookrightarrow X_0$ for all $X \in \text{ante } Q$. Hence, $X \vdash X_0$ for all $X \in \text{ante } Q$. By Proposition 3.4.2, we have $T_0 = T(X_0)$. Now let $X \in \text{ante } Q$ be given. Then $X \hookrightarrow X_0$, so by Proposition 3.4.1, we may choose $\psi \in \mathcal{F}$ such that $X \equiv X_0 \cup \{\psi\}$, and this gives $T(X) = T_0 + \psi$. \square

(P:basis-root) **Proposition 3.4.4.** *If Q is connected and \widehat{Q} is a basis for Q , then Q and \widehat{Q} have the same root.*

Proof. Let T_0 be the root of Q and \widehat{T}_0 be the root of \widehat{Q} . Then $T_0 = T(X_0)$ for some $X_0 \in \text{ante } Q$ and $\widehat{T}_0 = T(\widehat{X}_0)$ for some $\widehat{X}_0 \in \text{ante } \widehat{Q}$. Since $\widehat{Q} \subseteq Q$, we have $\text{ante } \widehat{Q} \subseteq \text{ante } Q$. Hence, $\widehat{X}_0 \in \text{ante } Q$, so by the definition of T_0 , we have $\widehat{X}_0 \vdash X_0$.

On the other hand, since $X_0 \in \text{ante } Q$ and \widehat{Q} is a basis for Q , we may choose $\widehat{X} \in \text{ante } \widehat{Q}$ and $S \subseteq \tau(\widehat{Q}; \widehat{X})$ such that $X_0 \equiv \widehat{X} \cup S$. Hence, $X_0 \vdash \widehat{X}$. But $\widehat{X} \in \text{ante } \widehat{Q}$, so by the definition of \widehat{T}_0 , we get $\widehat{X} \vdash \widehat{X}_0$. Thus, $X_0 \vdash \widehat{X}_0$, which shows that $X_0 \equiv \widehat{X}_0$, and therefore $T_0 = \widehat{T}_0$. \square

3.4.2 Connectivity and inductive inference

If P is admissible and strongly connected with root T_0 , and $\psi \in \mathcal{F}$ satisfies $T_0 + \psi \in \text{ante } P$, then we call ψ an *antecedent formula* of P . The set of all antecedent formulas of P is denoted by $\text{AF}(P)$. Note that

$$\text{ante } P = \{X \subseteq \mathcal{F}^{\text{IS}} \mid T(X) = T_0 + \psi \text{ for some } \psi \in \text{AF}(P)\}.$$

This follows from Proposition 3.4.3 and the rule of logical equivalence.

For this next result, recall the notation introduced in (3.2.6).

$\langle \text{P: simple-tau(P)} \rangle$ **Proposition 3.4.5.** *If P is entire and connected with root T_0 , then $\tau(P) = \tau(P; T_0)$.*

Proof. Let $\theta \in \tau(P)$. Then $P(\theta \mid X) = 1$ for all $X \in \text{ante } P$. In particular, $P(\theta \mid T_0) = 1$, so $\theta \in \tau(P; T_0)$. Conversely, suppose $P(\theta \mid T_0) = 1$ and let $X \in \text{ante } P$ be given. Since T_0 is the root of P , we have $X \vdash T_0$, so by deductive transitivity, $P(\theta \mid X) = 1$. Since X was arbitrary, $\theta \in \tau(P)$. \square

For this next result, recall the interval notation introduced just prior to Lemma 3.1.22.

$\langle \text{P: semi-cl-conn} \rangle$ **Proposition 3.4.6.** *Let P be semi-closed and connected with root T_0 . If $X \in \text{ante } P$, then $X \equiv T + \psi$ for some $T \in [T_0, \tau(P)]$ and some $\psi \in \mathcal{F}$. Moreover, ψ can be chosen so that $T_0 + \psi \in \text{ante } P$.*

Proof. Let \widehat{Q} be a basis for P . By Proposition 3.4.4, \widehat{Q} also has root T_0 . Let $X \in \text{ante } P$. Choose $\widehat{X} \in \text{ante } \widehat{Q}$ and $S \subseteq \tau(\widehat{Q}; \widehat{X})$ such that $X \equiv \widehat{X} \cup S$. By Proposition 3.4.3, we may choose $\psi \in \mathcal{F}$, such that $T(\widehat{X}) = T_0 + \psi$. Since $\widehat{Q} \subseteq P$, we have $\widehat{X} \in \text{ante } P$. By the rule of logical equivalence, $T_0 + \psi \in \text{ante } P$.

Now define $S' = \{\psi \rightarrow \theta \mid \theta \in S\}$ and $T = T_0 + S'$. Then $T_0 \subseteq T$ and, by Lemma 3.1.22, we have $X \equiv T(\widehat{X}) + S = T_0 + \psi + S = T_0 + \psi + S' = T + \psi$. Hence, it remains only to show that $T \subseteq \tau(P)$. By Proposition 3.4.5 and the fact that $T_0 \subseteq \tau(P; T_0)$, we need only show that $S' \subseteq \tau(P; T_0)$.

Let $\eta \in S'$. Choose $\theta \in S$ such that $\eta = \psi \rightarrow \theta$. Since $S \subseteq \tau(\widehat{Q}; \widehat{X})$, we have $\widehat{Q}(\theta \mid \widehat{X}) = 1$. Since $\widehat{Q} \subseteq P$, we have $P(\theta \mid \widehat{X}) = 1$. By the rule of logical equivalence, $P(\theta \mid T_0, \psi) = 1$. Hence, by the rule of material implication, it follows that $P(\eta \mid T_0) = P(\psi \rightarrow \theta \mid T_0) = 1$, showing that $\eta \in \tau(P; T_0)$. \square

3.4.3 Restrictions

For $Q \subseteq \mathcal{F}^{\text{IS}}$ and $\mathcal{X} \subseteq \text{ante } Q$, we define

$$Q \downarrow_{\mathcal{X}} = \{(X, \varphi, p) \in Q \mid X \hookrightarrow \mathcal{X}\}.$$

Note that if $Q \subseteq Q'$, then $Q \downarrow_{X_0} \subseteq Q' \downarrow_{X_0}$.

(T:restrict-inherit) **Theorem 3.4.7.** *Let $P \subseteq \mathcal{F}^{\text{IS}}$, $\mathcal{X} \subseteq \text{ante } P$, and define $P' = P \downarrow_{\mathcal{X}}$.*

(i) *If P is admissible, then P' is admissible.*

(ii) *If P is entire, then P' is entire.*

(iii) *If P is complete, then P' is complete.*

(iv) *If P is semi-closed, then P' is semi-closed.*

Proof. Assume P is admissible. Suppose $(X, \varphi, p) \in P'$, so that $X \hookrightarrow \mathcal{X}$ and $(X, \varphi, p) \in P$. Let $X' \equiv X$ and $\varphi' \equiv_X \varphi$. Since P is admissible, we have $(X', \varphi', p) \in P$, and since $X' \equiv X$, it follows that $X' \hookrightarrow \mathcal{X}$. Hence, $(X', \varphi', p) \in P'$. Now suppose $(X', \varphi', p') \in P'$. Then $(X', \varphi', p') \in P$, so the admissibility of P gives $p' = p$, and therefore P' is admissible.

Note that

(a) $X \in \text{ante } P'$ if and only if $X \in \text{ante } P$ and $X \hookrightarrow \mathcal{X}$,

(b) $P'(\varphi \mid X) = p$ if and only if $P(\varphi \mid X) = p$ and $X \hookrightarrow \mathcal{X}$, and

(c) $X \hookrightarrow \mathcal{X}$ implies $X \cup \{\varphi\} \hookrightarrow \mathcal{X}$.

Assume that P is entire. From (a) and (b), we easily see that P' satisfies (R2)–(R5) and (R7). From (a)–(c) we get that P' satisfies (R6), so P' is entire.

Assume P is complete. As above, (a) and (b) easily show that P' satisfies Definition 3.3.1, so that P' is complete.

Finally, assume P is semi-closed. Assume every completion of P' contains (X, φ, p) . Let \bar{P} be a completion of P . By (iii), $\bar{P} \downarrow_{\mathcal{X}}$ is complete. Since $P \subseteq \bar{P}$, we have $P' = P \downarrow_{\mathcal{X}} \subseteq \bar{P} \downarrow_{\mathcal{X}}$, so that $\bar{P} \downarrow_{\mathcal{X}}$ is a completion of P' . Hence, $\bar{P} \downarrow_{\mathcal{X}}(\varphi \mid X) = p$. By (b), we have $\bar{P}(\varphi \mid X) = p$ and $X \hookrightarrow \mathcal{X}$. Since \bar{P} was arbitrary and P is semi-closed, it follows that $P(\varphi \mid X) = p$. Since $X \hookrightarrow \mathcal{X}$, we get $P'(\varphi \mid X) = p$, and P' is semi-closed. \square

(C:restrict-inherit) **Corollary 3.4.8.** *If P is complete and $X_0 \in \text{ante } P$, then $P \downarrow_{X_0}$ is complete and strongly connected with root $T(X_0)$.*

Proof. This follows immediately from Theorem 3.4.7. \square

(C:restrict-ante) **Corollary 3.4.9.** *If P is entire and \bar{P} is a completion of P , then $\bar{P} \downarrow_{\text{ante } P}$ is also a completion of P .*

Proof. Let P be entire and let \bar{P} be a completion of P . Define $P' = \bar{P} \downarrow_{\text{ante } P}$. Theorem 3.4.7 implies that P' is also complete. Suppose $P(\varphi \mid X) = p$. Since $P \subseteq \bar{P}$, we have $\bar{P}(\varphi \mid X) = p$. Since $X \in \text{ante } P$, it follows that $P'(\varphi \mid X) = p$, showing that $P \subseteq P'$. \square

$\langle \text{C:ind-ext-str-conn} \rangle$ **Corollary 3.4.10.** *Suppose P is an entire set that has a completion. If $\bar{P}(\varphi \mid X) = p$ for all completions \bar{P} of P , then $X \leftrightarrow \text{ante } P$.*

Proof. Let P be an entire set that has a completion. Assume $\bar{P}(\varphi \mid X) = p$ for all completions \bar{P} of P . Choose a completion \bar{P} of P . By Corollary 3.4.9, the set $P' = \bar{P} \downarrow_{\text{ante } P}$ is also a completion of P . Hence, $P'(\varphi \mid X) = p$. By (b) above, we have $\bar{P}(\varphi \mid X) = p$ and $X \leftrightarrow \text{ante } P$. \square

$\langle \text{C:compl-restrict} \rangle$ **Corollary 3.4.11.** *Let P_0 be a pre-theory with root T_0 and let \bar{P}_0 be a completion of P_0 . Define $P'_0 = \bar{P}_0 \downarrow_{T_0}$. Then P'_0 is a completion of P_0 that is also a pre-theory with root T_0 .*

Proof. By Theorem 3.4.7, the set P'_0 is complete and strongly connected with root T_0 . Since complete sets are semi-closed, it follows that P'_0 is a pre-theory.

Suppose $P_0(\varphi \mid X) = p$. Since $P_0 \subseteq \bar{P}_0$, we have $\bar{P}_0(\varphi \mid X) = p$. Since P_0 is strongly connected with root T_0 , it follows that $X \leftrightarrow T_0$. Thus, $P'_0(\varphi \mid X) = p$, showing that $P_0 \subseteq P'_0$. \square

3.4.4 The lift of a pre-theory

Let P_0 be a pre-theory with root T_0 . By Proposition 3.2.13, the set $\tau(P_0)$ is a deductive theory. Since P_0 is admissible, we have $T_0 \in \text{ante } P_0$, so that by Proposition 3.4.5 and the rule of logical implication, $T_0 \subseteq \tau(P_0)$. Let $X \subseteq \mathcal{F}$. If $T(X) = T + \psi$ for some $T \in [T_0, \tau(P_0)]$ and some $\psi \in \text{AF}(P_0)$, then we call X a *generalized antecedent* of P_0 . The set of all such generalized antecedents is denoted by $\mathcal{GA}(P_0)$.

$\langle \text{P:lift-well-defined} \rangle$ **Proposition 3.4.12.** *Let P_0 be a pre-theory with root T_0 and let $X \in \mathcal{GA}(P_0)$. Suppose $T(X) = T + \psi = T' + \psi'$, where $T, T' \in [T_0, \tau(P_0)]$ and $\psi, \psi' \in \text{AF}(P_0)$. Then $P_0(\cdot \mid T_0, \psi) = P_0(\cdot \mid T_0, \psi')$.*

Proof. By symmetry, it suffices to show that $P_0(\varphi \mid T_0, \psi) = p$ implies $P_0(\varphi \mid T_0, \psi') = p$. Assume $P_0(\varphi \mid T_0, \psi) = p$. First note that $T + \psi = T' + \psi'$ implies $\psi' \in T + \psi \subseteq \tau(P_0) + \psi$, so that $\tau(P_0), \psi \vdash \psi'$. Likewise, $\tau(P_0), \psi' \vdash \psi$. Hence, $\psi \leftrightarrow \psi' \in \tau(P_0)$, giving $P_0(\psi \leftrightarrow \psi' \mid T_0) = 1$. Since $\psi' \in \text{AF}(P_0)$, we have $T_0 + \psi' \in \text{ante } P_0$, so by two applications of deductive transitivity, $P_0(\psi' \rightarrow \psi \mid T_0, \psi') = 1$. By the rule of logical implication, $P_0(\psi' \mid T_0, \psi') = 1$. Applying Proposition 3.2.16(iii) with $\zeta = \top$ gives $P_0(\psi \mid T_0, \psi') = 1$. A similar argument shows that $P_0(\psi' \mid T_0, \psi) = 1$.

Now, by Proposition 3.2.11, we have $P_0(\varphi \wedge \psi' \mid T_0, \psi) = p$. By the multiplication rule, $P_0(\varphi \mid T_0, \psi, \psi') = p$. A second application of the multiplication rule then gives $P_0(\varphi \wedge \psi \mid T_0, \psi') = p$. By Proposition 3.2.9, it follows that $P_0(\varphi \mid T_0, \psi') = p$. \square

Let P_0 be a pre-theory with root T_0 . For each $X \in \mathcal{GA}(P_0)$, choose $T^X \in [T_0, \tau(P_0)]$ and $\psi_X \in \text{AF}(P_0)$ such that $T(X) = T^X + \psi_X$. Define

$$\mathbf{L}(P_0) = \{(X, \varphi, p) \mid X \in \mathcal{GA}(P_0), P_0(\varphi \mid T_0, \psi_X) = p\}.$$

By Proposition 3.4.12, the definition of $\mathbf{L}(P_0)$ does not depend on the choices of T^X and ψ_X . We call $\mathbf{L}(P_0)$ the *lift* of P_0 , and may sometimes denote it by \mathbf{L}_{P_0} .

$\langle \text{P: any-psi} \rangle$ **Proposition 3.4.13.** *Let P_0 be a pre-theory with root T_0 . If $X \in \text{ante } \mathbf{L}(P_0)$ and $T(X) = T + \psi$, where $T \in [T_0, \tau(P_0)]$ and $\psi \in \mathcal{F}$, then $\psi \in \text{AF}(P_0)$.*

Proof. Since $X \in \text{ante } \mathbf{L}(P_0)$, we may choose φ and p such that $(X, \varphi, p) \in \mathbf{L}(P_0)$. Then $T(X) = T^X + \psi_X$ and $P_0(\varphi \mid T_0, \psi_X) = p$. As in the proof of Proposition 3.4.12, since $T, T^X \subseteq \tau(P_0)$ and $T + \psi = T^X + \psi_X$, we have $\psi \leftrightarrow \psi_X \in \tau(P_0)$. Thus, $P_0(\psi \leftrightarrow \psi_X \mid T_0) = 1$. Since P_0 is semi-closed, Proposition 3.3.3 gives $P_0(\varphi \mid T_0, \psi) = p$, showing that $T_0 + \psi \in \text{ante } P_0$, and therefore $\psi \in \text{AF}(P_0)$. \square

$\langle \text{P: lift-extends} \rangle$ **Proposition 3.4.14.** *If P_0 is a pre-theory, then $P_0 \subseteq \mathbf{L}(P_0)$.*

Proof. Let T_0 be the root of P_0 . Suppose $P_0(\varphi \mid X) = p$. Since P_0 is strongly connected, we may choose $\varphi \in \mathcal{F}$ such that $T(X) = T_0 + \psi$. By the rule of logical equivalence, $P_0(\varphi \mid T_0, \psi) = p$. Thus, $\psi \in \text{AF}(P_0)$ and, since $T_0 \in [T_0, \tau(P_0)]$, we have $X \in \mathcal{GA}(P_0)$. Therefore, $(X, \varphi, p) \in \mathbf{L}(P_0)$. \square

$\langle \text{P: lift-unique} \rangle$ **Proposition 3.4.15.** *If P_0 is a pre-theory with root T_0 , then $\mathbf{L}(P_0) \downarrow_{T_0} = P_0$.*

Proof. From Proposition 3.4.14, it follows that $P_0 = P_0 \downarrow_{T_0} \subseteq \mathbf{L}(P_0) \downarrow_{T_0}$. For the reverse, suppose $(X, \varphi, p) \in \mathbf{L}(P_0) \downarrow_{T_0}$. Then $X \in \mathcal{GA}(P_0)$, $T(X) = T^X + \psi_X$, $P_0(\varphi \mid T_0, \psi_X) = p$, and $X \hookrightarrow T_0$. By Proposition 3.4.1, we may write $T(X) = T_0 + \psi$ for some $\psi \in \mathcal{F}$. By Proposition 3.4.13, we know that $\psi \in \text{AF}(P_0)$. Therefore, by Proposition 3.4.12, we have $P_0(\varphi \mid T_0, \psi) = p$. By the rule of logical equivalence, $P_0(\varphi \mid X) = p$. \square

$\langle \text{L: lift-admissible} \rangle$ **Lemma 3.4.16.** *Let P_0 be a pre-theory. If $X \in \text{ante } \mathbf{L}(P_0)$ and $X \vdash \varphi$, then $(X, \varphi, 1) \in \mathbf{L}(P_0)$.*

Proof. Choose λ and p such that $(X, \lambda, p) \in \mathbf{L}(P_0)$. Then $T(X) = T^X + \psi_X$ and $P_0(\lambda \mid T_0, \psi_X) = p$. In particular, $T_0 + \psi_X \in \text{ante } P_0$.

Since $X \vdash \varphi$, we have $\varphi \in T(X) = T^X + \psi_X$, meaning $T^X, \psi_X \vdash \varphi$. Choose countable $\Phi \subseteq T^X$ such that $\Phi, \psi_X \vdash \varphi$. Since $\Phi \subseteq T^X \subseteq \tau(P_0)$, it follows that $P_0(\theta \mid T_0) = 1$ for all $\theta \in \Phi$. By Corollary 3.2.12, if $\zeta = \bigwedge \Phi$, then $P_0(\zeta \mid T_0) = 1$. Since $T_0 + \psi_X \in \text{ante } P_0$, deductive transitivity gives $P_0(\zeta \mid T_0, \psi_X) = 1$. By the rule of logical implication, $P_0(\psi_X \mid T_0, \psi_X) = 1$. Hence, by Proposition 3.2.11, we obtain $P_0(\zeta \wedge \psi_X \mid T_0, \psi_X) = 1$. But $\zeta \wedge \psi_X \vdash \varphi$, so by deductive transitivity, $P_0(\varphi \mid T_0, \psi_X) = 1$, which gives $(X, \varphi, 1) \in \mathbf{L}(P_0)$. \square

$\langle \text{P: lift-admissible} \rangle$ **Proposition 3.4.17.** *If P_0 is a pre-theory, then $\mathbf{L}(P_0)$ is admissible.*

Proof. Suppose $(X, \varphi, p) \in \mathbf{L}(P_0)$, $X' \equiv X$, and $\varphi' \equiv_X \varphi$. Then $T(X) = T^X + \psi_X$ and $P_0(\varphi \mid T_0, \psi_X) = p$. Since $\varphi' \equiv_X \varphi$, we have $X \vdash \varphi' \leftrightarrow \varphi$. From Lemma 3.4.16, it follows that $(X, \varphi' \leftrightarrow \varphi, 1) \in \mathbf{L}(P_0)$, which gives $P_0(\varphi' \leftrightarrow \varphi \mid T_0, \psi_X) = 1$. By Proposition 3.2.14, we have $P_0(\varphi' \mid T_0, \psi_X) = p$. But $X' \equiv X$, so $T(X') = T(X) = T^X + \psi_X$. Thus, $(X', \varphi', p) \in \mathbf{L}(P_0)$. Finally, if $(X', \varphi', p') \in \mathbf{L}(P_0)$, then $P_0(\varphi' \mid T_0, \psi_X) = p'$. But P_0 is admissible, so $p' = p$. \square

Let P_0 be a pre-theory with root T_0 and let $P = \mathbf{L}(P_0)$. By Proposition 3.4.17, we may now use the notation $P(\varphi \mid X) = p$ instead of $(X, \varphi, p) \in P$. Note that by Proposition 3.4.13, if $X \in \text{ante } P$, $T \in [T_0, \tau(P_0)]$, $\psi \in \mathcal{F}$, and $T(X) = T + \psi$, then $P(\cdot \mid X) = P_0(\cdot \mid T_0, \psi)$.

$\langle \text{P:nested-lifts} \rangle$ **Proposition 3.4.18.** *Let P_0, P'_0 be pre-theories with a common root. If $P_0 \subseteq P'_0$, then $\mathbf{L}(P_0) \subseteq \mathbf{L}(P'_0)$.*

Proof. Let T_0 be the common root of P_0 and P'_0 . Assume $P_0 \subseteq P'_0$. Let $P = \mathbf{L}(P_0)$ and $P' = \mathbf{L}(P'_0)$. Suppose $P(\varphi \mid X) = p$. Then $T(X) = T^X + \psi_X$ and $P_0(\varphi \mid T_0, \psi_X) = p$. Since $P_0 \subseteq P'_0$, we have $\tau(P_0) \subseteq \tau(P'_0)$, and therefore $[T_0, \tau(P_0)] \subseteq [T_0, \tau(P'_0)]$. From $P_0 \subseteq P'_0$, it also follows that $\text{AF}(P_0) \subseteq \text{AF}(P'_0)$. Hence, $T^X \in [T_0, \tau(P'_0)]$ and $\psi_X \in \text{AF}(P'_0)$. This shows that $X \in \mathcal{GA}(P'_0)$ and $P'(\varphi \mid X) = P'_0(\varphi \mid T_0, \psi_X)$. But $P_0 \subseteq P'_0$, so $P'_0(\varphi \mid T_0, \psi_X) = p$. \square

3.4.5 Identifying lifts with inductive theories

We are now ready to prove Theorem 3.3.4. We will do this by showing that $\mathbf{L}(P_0)$ is the smallest closed extension of P_0 . The proof of Theorem 3.3.4 is broken into five pieces for greater readability.

$\langle \text{P:lift-entire} \rangle$ **Proposition 3.4.19.** *If P_0 is a pre-theory, then $\mathbf{L}(P_0)$ is entire.*

Proof. Let P_0 be a pre-theory with root T_0 and let $P = \mathbf{L}(P_0)$. By Proposition 3.4.17, P is admissible, and by Lemma 3.4.16, P satisfies the rule of logical implication. It therefore remains only to verify (R3)–(R7).

We begin with the rule of material implication. Suppose $X \in \text{ante } P$ and $P(\varphi \mid X, \psi) = 1$. Then $T(X) = T^X + \psi_X$ and $P(\cdot \mid X) = P_0(\cdot \mid T_0, \psi_X)$. Thus, $T(X \cup \{\psi\}) = T^X + \psi_X \wedge \psi$, so by Proposition 3.4.13, we have $P_0(\varphi \mid T_0, \psi_X \wedge \psi) = 1$. Let $T' = T_0 + \psi_X$, so that $T' + \psi = T_0 + \psi_X \wedge \psi$. Then $T' \in \text{ante } P_0$ and $P_0(\varphi \mid T', \psi) = 1$. By the rule of material implication for P_0 , it follows that $P_0(\psi \rightarrow \varphi \mid T') = 1$, which implies $P(\psi \rightarrow \varphi \mid X) = P_0(\psi \rightarrow \varphi \mid T_0, \psi_X) = P_0(\psi \rightarrow \varphi \mid T') = 1$.

We next verify deductive transitivity. Suppose $P(\varphi \mid X) = 1$ and $\varphi \vdash \psi$. Then $P_0(\varphi \mid T_0, \psi_X) = 1$. By the deductive transitivity for P_0 , we have $P_0(\psi \mid T_0, \psi_X) = 1$, which gives $P(\psi \mid X) = 1$. Now suppose $X' \in \text{ante } P$, $X' \vdash X$, and $P(\varphi \mid X) = 1$. Then $P_0(\varphi \mid T_0, \psi_X) = 1$ and, since $X' \in \text{ante } P$, we get $T(X') = T^{X'} + \psi_{X'}$ and $P(\cdot \mid X') = P_0(\cdot \mid T_0, \psi_{X'})$. Since $X' \vdash X$, we have $T^{X'}, \psi_{X'} \vdash \psi_X$. Choose countable $\Phi \subseteq T^{X'}$ such that $\Phi, \psi_{X'} \vdash \psi_X$ and let $\zeta = \bigwedge \Phi$. Then $\Phi \subseteq T^{X'} \subseteq \tau(P_0)$, so Corollary 3.2.12 implies $P_0(\zeta \mid T_0) = 1$.

Since $T_0 + \psi_{X'} \in \text{ante } P_0$, deductive transitivity for P_0 gives $P_0(\zeta \mid T_0, \psi_{X'}) = 1$. By Lemma 3.2.10, we have $T_0 + \psi_{X'} + \zeta \in \text{ante } P_0$. Since $T_0 + \psi_{X'} + \zeta \vdash T_0, \psi_X$, we may again apply deductive transitivity for P_0 to obtain $P_0(\varphi \mid T_0, \psi_{X'}, \zeta) = 1$. By the multiplication rule, $P_0(\zeta \wedge \varphi \mid T_0, \psi_{X'}) = 1$. By Proposition 3.2.9, we obtain $P_0(\varphi \mid T_0, \psi_{X'}) = P_0(\zeta \wedge \varphi \mid T_0, \psi_{X'}) = 1$, and this shows that $P(\varphi \mid X') = 1$.

To show that P satisfies the addition rule, suppose $X \vdash \neg(\varphi \wedge \psi)$ and assume two of the probabilities in (3.2.1) exist. Then $X \in \text{ante } P$, so that $T(X) = T^X + \psi_X$ and $P(\cdot \mid X) = P_0(\cdot \mid T_0, \psi_X)$. Thus, two of the probabilities in the following equation exist:

$$P_0(\varphi \vee \psi \mid T_0, \psi_X) = P_0(\varphi \mid T_0, \psi_X) + P_0(\psi \mid T_0, \psi_X).$$

By the addition rule for P_0 , so does the third, and the above equation holds. Hence, all three probabilities in (3.2.1) exist and (3.2.1) holds. The proof that P satisfies the multiplication rule is similar.

Finally, suppose $P(\varphi_n \mid X)$ exists and $X, \varphi_n \vdash \varphi_{n+1}$ for all n . We want to show that (3.2.3) holds. Since $X \in \text{ante } P$, we have $T(X) = T^X + \psi_X$ and $P(\cdot \mid X) = P_0(\cdot \mid T_0, \psi)$. Thus, $P_0(\varphi_n \mid T_0, \psi)$ exists for all n .

First assume $P_0(\varphi_n \mid T_0, \psi) = 0$ for all n . Let $\psi_n = \bigvee_1^n \varphi_j$. Then $P_0(\psi_n \mid T_0, \psi) = 0$ by Proposition 3.2.11. Therefore, $P_0(\bigvee_n \psi_n \mid T_0, \psi) = 0$ by the continuity rule. But $\bigvee_n \psi_n \equiv \bigvee_n \varphi_n$, so $P_0(\bigvee_n \varphi_n \mid T_0, \psi) = 0$, which implies $P(\bigvee_n \varphi_n \mid X) = 0$, and (3.2.3) holds in this case.

Now assume there exists n_0 such that $P_0(\varphi_{n_0} \mid T_0, \psi) > 0$. By Remark 3.2.8, we have that P satisfies Proposition 3.2.5. Thus, $P(\varphi_n \mid X) \leq P(\varphi_{n+1} \mid X)$, which implies $P_0(\varphi_n \mid T_0, \psi) \leq P_0(\varphi_{n+1} \mid T_0, \psi)$. Hence, $P_0(\varphi_n \mid T_0, \psi) > 0$ for all $n \geq n_0$. Since $X \vdash \varphi_n \rightarrow \varphi_{n+1}$, we have $P(\varphi_n \rightarrow \varphi_{n+1} \mid X) = 1$, which gives $P_0(\varphi_n \rightarrow \varphi_{n+1} \mid T_0, \psi) = 1$. From Proposition 3.2.15, it follows that $P_0(\varphi_{n+1} \mid T_0, \psi, \varphi_n) = 1$. Therefore, by Theorem 3.2.23, we have $P_0(\bigvee_{n_0}^\infty \varphi_n \mid T_0, \psi) = \lim_n P_0(\varphi_n \mid T_0, \psi)$, which implies $P(\bigvee_{n_0}^\infty \varphi_n \mid X) = \lim_n P(\varphi_n \mid X)$. But $\bigvee_{n_0}^\infty \varphi_n \equiv_X \bigvee_n \varphi_n$, so (3.2.3) holds in this case as well. \square

(P:lift-complete) **Proposition 3.4.20.** *If P_0 is a complete pre-theory, then $\mathbf{L}(P_0)$ is complete.*

Proof. Let P_0 be a complete pre-theory with root T_0 and let $P = \mathbf{L}(P_0)$. Proposition 3.4.19 implies P is entire. Suppose $P(\varphi \mid X) = p$ and $P(\psi \mid X) = q$. Then $T(X) = T^X + \psi_X$, $P_0(\varphi \mid T_0, \psi_X) = p$, and $P_0(\psi \mid T_0, \psi_X) = q$. Since P_0 is complete, $P_0(\varphi \wedge \psi \mid T_0, \psi_X) = r$ for some r . Thus, $P(\varphi \wedge \psi \mid X) = r$, and P satisfies Definition 3.3.1(i).

Now suppose $X \in \text{ante } P$ and $X \cup \{\varphi\} \in \text{ante } P$. Then $T(X) = T^X + \psi_X$, so that $T(X \cup \{\varphi\}) = T(X) + \varphi = T^X + \psi_X + \varphi = T^X + \psi_X \wedge \varphi$. By Proposition 3.4.13, we have $\psi_X \wedge \varphi \in \text{AF}(P_0)$. Thus, $T_0 + \psi_X \wedge \varphi \in \text{ante } P_0$. Since $T_0 + \psi_X \wedge \varphi \equiv (T_0 + \psi_X) \cup \{\varphi\}$, the rule of logical equivalence gives $(T_0 + \psi_X) \cup \{\varphi\} \in \text{ante } P_0$. But $T_0 + \psi_X \in \text{ante } P_0$ and P_0 is complete, so $P_0(\varphi \mid T_0, \psi_X)$ exists. Hence, $P(\varphi \mid X) = P_0(\varphi \mid T_0, \psi_X)$ exists, so that P satisfies Definition 3.3.1(ii). \square

$\langle \text{P:lift-semi-closed} \rangle$ **Proposition 3.4.21.** *If P_0 is a pre-theory, then $\mathbf{L}(P_0)$ is semi-closed.*

Proof. Let P_0 be a pre-theory with root T_0 and let $P = \mathbf{L}(P_0)$. Proposition 3.4.19 implies P is entire. Suppose $\bar{P}(\varphi \mid X) = p$ whenever \bar{P} is a completion of P . Since P_0 is a pre-theory, it is semi-closed. It therefore has a completion. Let \bar{P}_0 be a completion of P_0 . Corollary 3.4.11 implies $P'_0 = \bar{P}_0 \downarrow_{T_0}$ is a completion of P_0 that is also a pre-theory with root T_0 . Let $P' = \mathbf{L}(P'_0)$. Then Proposition 3.4.20 implies P' is complete and Proposition 3.4.18 implies $P \subseteq P'$. Hence, P' is a completion of P , so by supposition, $P'(\varphi \mid X) = p$.

By Corollary 3.4.10, we have $X \leftrightarrow \text{ante } P$. Choose $X' \in \text{ante } P$ and $\psi' \in \mathcal{F}$ such that $X \equiv X' \cup \{\psi'\}$. Then $T(X) = T(X') + \psi' = T^{X'} + \psi_{X'} + \psi' = T^{X'} + \psi$, where $\psi = \psi_{X'} \wedge \psi'$. From Proposition 3.4.13, it follows that $\psi \in \text{AF}(P_0)$. As in the proof of Proposition 3.4.18, we have $[T_0, \tau(P_0)] \subseteq [T_0, \tau(P'_0)]$ and $\text{AF}(P_0) \subseteq \text{AF}(P'_0)$. Hence, $T^{X'} \in [T_0, \tau(P'_0)]$ and $\psi \in \text{AF}(P'_0)$, which implies $P'_0(\varphi \mid T_0, \psi) = p$. But $P'_0 \subseteq \bar{P}_0$, so $\bar{P}_0(\varphi \mid T_0, \psi) = p$. Since \bar{P}_0 was arbitrary and P_0 is semi-closed, the rule of inductive extension gives $P_0(\varphi \mid T_0, \psi) = p$. Hence, $P(\varphi \mid X) = P(\varphi \mid T^{X'}, \psi) = p$, which shows that P satisfies the rule of induction extension and is therefore semi-closed. \square

$\langle \text{P:lift-closed} \rangle$ **Proposition 3.4.22.** *If P_0 is a pre-theory, then $\mathbf{L}(P_0)$ is a closed extension of P_0 .*

Proof. Let P_0 be a pre-theory with root T_0 and let $P = \mathbf{L}(P_0)$. Proposition 3.4.14 implies P is an extension of P_0 , and Proposition 3.4.21 implies P is semi-closed. Let $S \subseteq \mathcal{F}$ be nonempty and assume $P(\theta \mid X) = 1$ for all $\theta \in S$. Then $X \in \text{ante } P$ and $P_0(\theta \mid T_0, \psi_X) = 1$ for all $\theta \in S$. By the rule of material implication, $P_0(\psi_X \rightarrow \theta \mid T_0) = 1$ for all $\theta \in S$. Hence, if we define $S' = \{\psi_X \rightarrow \theta \mid \theta \in S\}$ and $T' = T^X + S'$, then Proposition 3.4.5 implies $S' \subseteq \tau(P_0)$, so that $T' \in [T_0, \tau(P_0)]$. By Lemma 3.1.22, we have $X \cup S \equiv T^X + \psi_X + S = T^X + \psi_X + S' = T' + \psi_X$. It therefore follows that $X \cup S \in \text{ante } P$ and $P(\varphi \mid X \cup S) = p$ if and only if $P_0(\varphi \mid T_0, \psi_X) = p$, which holds if and only if $P(\varphi \mid X) = p$. This shows that P satisfies the rule of deductive extension and is therefore closed. \square

Proof of Theorem 3.3.4. Let P_0 be a pre-theory with root T_0 and let $P = \mathbf{L}(P_0)$. Proposition 3.4.22 shows that P is a closed extension of P_0 . Let P' be another closed extension of P_0 and suppose $P(\varphi \mid X) = p$. Then $P_0(\varphi \mid T_0, \psi_X) = p$, which implies $P'(\varphi \mid T_0, \psi_X) = p$. Since $T^X \subseteq \tau(P_0)$, we have $P_0(\theta \mid T_0) = 1$ for all $\theta \in T^X$. Deductive transitivity gives $P_0(\theta \mid T_0, \psi_X) = 1$ for all $\theta \in T^X$, which implies $P'(\theta \mid T_0, \psi_X) = 1$ for all $\theta \in T^X$. Since P' is closed, the rule of deductive extension gives $P'(\varphi \mid T_0, \psi_X, T^X) = 1$. But $T_0 \subseteq T^X$, so $T_0 + \psi_X + T^X = T^X + \psi_X \equiv X$, and the rule of logical equivalence gives $P'(\varphi \mid X) = p$. \square

3.4.6 Characterizing inductive theories

Recall the notation $\mathbf{P}(P_0)$, established in Section 3.3.4. The proof of Theorem 3.3.4 shows that $\mathbf{P}(P_0) = \mathbf{L}(P_0)$, so that $\mathbf{P}(P_0)$ is simply the lift of P_0 . In particular, this gives us an explicit construction of $\mathbf{P}(P_0)$ from P_0 . We will use this explicit construction to prove Corollary 3.3.5 and Theorem 3.3.6. Henceforth, we will drop the notation $\mathbf{L}(P_0)$, and write only $\mathbf{P}(P_0)$ instead.

Proof of Corollary 3.3.5. Let P_0, P'_0 be pre-theories with roots T_0, T'_0 , respectively, and let $P = \mathbf{P}(P_0)$ and $P' = \mathbf{P}(P'_0)$.

First assume $P = P'$. Then $T_0 \in \text{ante } P_0 \subseteq \text{ante } P = \text{ante } P'$. We may therefore choose $T' \in [T'_0, \tau(P'_0)]$ and $\psi' \in \text{AF}(P'_0)$ such that $T_0 = T' + \psi'$. Then $T'_0 \subseteq T' \subseteq T' + \psi' = T_0$. By reversing the roles of P_0 and P'_0 , we also have $T_0 \subseteq T'_0$, showing that $T_0 = T'_0$. Thus, by Proposition 3.4.15, we have $P_0 = P \downarrow_{T_0} = P' \downarrow_{T'_0} = P'_0$.

Now assume $P_0 \subset P'_0 \subseteq P$. Then P is a closed extension of P'_0 , so Theorem 3.3.4 implies $P' \subseteq P$. On the other hand, P' is a closed extension of P'_0 and $P_0 \subset P'_0$, so P' is a closed extension of P_0 , giving $P \subseteq P'$. Thus, $P = P'$, so by the above $P_0 = P'_0$, a contradiction. \square

$\langle \mathbf{P}:\text{lift-connected} \rangle$ **Proposition 3.4.23.** *If P_0 is a pre-theory, then $\mathbf{P}(P_0)$ is connected.*

Proof. Let P_0 be a pre-theory with root T_0 and let $P = \mathbf{P}(P_0)$. We will show that P_0 is a basis for P . First note that P_0 is strongly connected and $P_0 \subseteq P$. Now let $X \in \text{ante } P$ be given. Define $\widehat{X} = T_0 + \psi_X$ and $S = T^X$. Then $\widehat{X} \in \text{ante } P_0$. Since $T^X \subseteq \tau(P_0)$, we have $P_0(\theta \mid T_0) = 1$ for all $\theta \in S$. By deductive transitivity, $P_0(\theta \mid \widehat{X}) = 1$ for all $\theta \in S$. Hence, $S \subseteq \tau(P_0, \widehat{X})$. Lastly, from $T_0 \subseteq T^X$, it follows that $X \equiv T^X + \psi_X = T_0 + T^X + \psi_X \equiv \widehat{X} \cup S$. \square

Proof of Theorem 3.3.6. By Theorem 3.3.4 and Proposition 3.4.23, we have (i) implies (ii). For the converse, let P be closed and connected. Let T_0 be the root of P and define $P_0 = P \downarrow_{T_0}$. By Theorem 3.4.7, the set P_0 is semi-closed. By construction, P_0 is strongly connected with root T_0 . Thus, P_0 is a pre-theory. We will show that $P = \mathbf{P}(P_0)$.

For notational simplicity, let $P' = \mathbf{P}(P_0)$. Since P is a closed extension of P_0 , Theorem 3.3.4 gives $P' \subseteq P$. Suppose $P(\varphi \mid X) = p$. By Proposition 3.4.6, we may choose $T \in [T_0, \tau(P)]$ and $\psi \in \mathcal{F}$ such that $X \equiv T + \psi$ and $T_0 + \psi \in \text{ante } P$. Since $T \subseteq \tau(P)$, we have $P(\theta \mid T_0) = 1$ for all $\theta \in T$. By deductive transitivity, $P(\theta \mid T_0, \psi) = 1$ for all $\theta \in T$. Since P is closed, the rule of deductive extension implies $P(\cdot \mid T_0, \psi) = P(\cdot \mid T_0, \psi, T)$. But $T_0 + \psi + T = T + \psi \equiv X$. Hence, $P(\cdot \mid T_0, \psi) = P(\cdot \mid X)$, and it therefore follows that $P(\varphi \mid T_0, \psi) = p$. Since $P_0 = P \downarrow_{T_0}$, this gives $P_0(\varphi \mid T_0, \psi) = p$. Finally, since P' is the lift of P_0 , we have $P'(\varphi \mid X) = P'(\varphi \mid T, \psi) = P_0(\varphi \mid T_0, \psi) = p$. \square

$\langle \mathbf{R}:\text{ante-cao} \rangle$ **Remark 3.4.24.** Note that by Theorem 3.3.6 and Proposition 3.4.15, if P is an inductive theory, then $P = \mathbf{P}(P_0)$, where $P_0 = P \downarrow_{T_0}$. Also note that by Proposition 3.4.5, we have $\tau(P) = \tau(P_0)$. Hence, every $X \in \text{ante } P$ satisfies $X \leftrightarrow [T_0, \tau(P)]$.

3.5 Generating inductive theories

$\langle \mathbf{S} : \text{ind-cond} \rangle$ In the first half of this section, we prove Theorem 3.3.7. One might think that we could do this the same way we would do it for the deductive calculus. Namely, we could try to prove that (i) the intersection of inductive theories is an inductive theory; therefore, (ii) the intersection of all inductive theories that contain Q is the smallest inductive theory containing Q . Unfortunately, as we will see in Section 4.4.1 (see Remark 4.4.3), it turns out that (i) is false. (Surprisingly, though, (ii) is still true, provided we pay special attention to the root of Q .) We therefore have to find a different proof method.

As it turns out, the intersection of pre-theories (with a common root) is a pre-theory. This will be the key to our proof, but it will require several preliminary definitions and results. When we are done proving Theorem 3.3.7, we present a partial converse to the rule of logical implication (see Theorem 3.5.6).

In the second half of this section, we generalize inductive derivability to something that we call “inductive conditions”. This allows us to reason with more general statements, such as $Q(\varphi \mid X) > 1/2$.

3.5.1 Strongly connected equivalence

Our first preliminary on the way to the proof of Theorem 3.3.7 shows that, when extending a consistent set to an inductive theory, we are never required to change roots.

$\langle \mathbf{P} : \text{cons-match-root} \rangle$ **Proposition 3.5.1.** *Every consistent set can be extended to an inductive theory with the same root. More specifically, let Q be consistent with root T_0 and let P be an inductive theory with $Q \subseteq P$. Then $T_0 \in \text{ante } P$ and $P \downarrow_{T_0}$ is a pre-theory with root T_0 that satisfies $Q \subseteq \mathbf{P}(P \downarrow_{T_0}) \subseteq P$.*

Proof. Let Q be consistent with root T_0 and let P be an inductive theory with $Q \subseteq P$. Since T_0 is the root of Q , we have $T_0 = T(X_0)$ for some $X_0 \in \text{ante } Q$. But $Q \subseteq P$, so $X_0 \in \text{ante } P$ and, by the rule of logical equivalence, $T_0 \in \text{ante } P$. Let $P'_0 = P \downarrow_{T_0}$. Then P'_0 is strongly connected with root T_0 and, by Theorem 3.4.7, the set P'_0 is semi-closed. Therefore, P'_0 is a pre-theory, and we may define $P' = \mathbf{P}(P'_0)$. Note that P' is an inductive theory with root T_0 . Since $P'_0 \subseteq P$ and P is an inductive theory, Theorem 3.3.4 implies $P' \subseteq P$. It remains only to show that $Q \subseteq P'$.

Suppose $Q(\varphi \mid X) = p$. We want to show that $P'(\varphi \mid X) = p$. Since P' is the lift of P'_0 , we must find $T \in [T_0, \tau(P'_0)]$ and $\psi \in \mathcal{F}$ such that $X \equiv T + \psi$ and $P'_0(\varphi \mid T_0, \psi) = p$. Let \widehat{Q} be a basis for Q . Choose $\widehat{X} \in \text{ante } \widehat{Q}$ and $S \subseteq \tau(\widehat{Q}; \widehat{X})$ such that that $X \equiv \widehat{X} \cup S$. Proposition 3.4.4 implies that T_0 is the root of \widehat{Q} . Hence, we may choose $\psi \in \mathcal{F}$ such that $\widehat{X} \equiv T_0 + \psi$. Define $S' = \{\psi \rightarrow \theta \mid \theta \in S\}$ and $T = T_0 + S'$. By Lemma 3.1.22, we have $X \equiv \widehat{X} \cup S \equiv T_0 + \psi + S = T_0 + \psi + S' = T + \psi$, so it suffices to show $S' \subseteq \tau(P'_0)$ and $P'_0(\varphi \mid T_0, \psi) = p$.

By Proposition 3.4.5, in order to show that $S' \subseteq \tau(P'_0)$, we must show that $P'_0(\psi \rightarrow \theta \mid T_0) = 1$ for all $\theta \in S$. Let $\theta \in S$. Then $\widehat{Q}(\theta \mid \widehat{X}) = 1$. But $\widehat{Q} \subseteq Q \subseteq P$, so $P(\theta \mid \widehat{X}) = P(\theta \mid T_0, \psi) = 1$. By the rule of material implication, $P(\psi \rightarrow \theta \mid T_0) = 1$, which implies $P'_0(\psi \rightarrow \theta \mid T_0) = 1$.

Finally, since $\widehat{Q} \subseteq Q \subseteq P$, we have $\tau(\widehat{Q}; \widehat{X}) \subseteq \tau(P; \widehat{X})$. Hence, $S \subseteq \tau(P; \widehat{X})$. Since P is closed, deductive extension implies $P(\cdot \mid \widehat{X}) = P(\cdot \mid \widehat{X} \cup S)$. But $\widehat{X} \cup S \equiv X$. Also, $Q \subseteq P$, so we have $P(\varphi \mid X) = p$. Using the rule of logical equivalence, we have $P(\varphi \mid T_0, \psi) = P(\varphi \mid \widehat{X}) = P(\varphi \mid X) = p$. Therefore, since $P'_0 = P \downarrow_{T_0}$, we obtain $P'_0(\varphi \mid T_0, \psi) = p$. \square

Our next result says that connected sets are, in a certain sense, logically unnecessary. It is enough to only consider strongly connected sets. To make this precise, we first define what it means for two subsets of \mathcal{F}^{IS} to be logically equivalent.

Let $Q, Q' \subseteq \mathcal{F}^{\text{IS}}$ be connected. We say that Q and Q' are *equivalent*, written $Q \equiv Q'$, if, for all inductive theories P , we have $Q \subseteq P$ if and only if $Q' \subseteq P$. After proving Theorem 3.3.7, we will be able to speak of the inductive theory generated by Q . At that point, we will have the much more natural characterization of equivalence given in Proposition 3.5.5.

(T:str-conn-enough) **Theorem 3.5.2.** *If Q is connected, then there exists strongly connected Q_0 with the same root as Q such that $Q_0 \equiv Q$.*

Proof. Let Q be connected with root T_0 and let \widehat{Q} be a basis for Q . For each $X \in \text{ante } Q$, choose $\widehat{X} \in \text{ante } \widehat{Q}$, $S_X \subseteq \tau(\widehat{Q}; \widehat{X})$, and $\psi_X \in \mathcal{F}$ such that $X \equiv \widehat{X} \cup S_X$ and $\widehat{X} \equiv T_0 + \psi_X$. For $X \in \text{ante } \widehat{Q} \subseteq \text{ante } Q$, assume we have chosen $\widehat{X} = X$ and $S_X = \emptyset$. Define

$$Q_0 = \{(T_0 + \psi_X, \varphi, p) \mid (X, \varphi, p) \in Q\}.$$

Note that if $\widehat{X} \in \text{ante } \widehat{Q} \subseteq \text{ante } Q$, then $S_{\widehat{X}} = \emptyset$ and $\widehat{X} \equiv T_0 + \psi_{\widehat{X}}$. Thus, if $(\widehat{X}, \varphi, p) \in \widehat{Q} \subseteq Q$, then $(T_0 + \psi_{\widehat{X}}, \varphi, p) \in Q_0$.

Let P be an inductive theory. Assume $Q_0 \subseteq P$. Let $(X, \varphi, p) \in Q$. Then $(T_0 + \psi_X, \varphi, p) \in Q_0$, which implies $P(\varphi \mid T_0, \psi_X) = p$. Let $\theta \in S_X \subseteq \tau(\widehat{Q}; \widehat{X})$. Then $(\widehat{X}, \theta, 1) \in \widehat{Q}$, so that $(T_0 + \psi_{\widehat{X}}, \theta, 1) \in Q_0$, which implies $P(\theta \mid T_0, \psi_{\widehat{X}}) = 1$. But $T_0 + \psi_{\widehat{X}} \equiv \widehat{X} \equiv T_0 + \psi_X$, so by the rule of logical equivalence, $P(\theta \mid T_0, \psi_X) = 1$. Since θ was arbitrary, deductive extension gives $P(\varphi \mid T_0, \psi_X, S_X) = p$. But $T_0 + \psi_X + S_X \equiv \widehat{X} \cup S_X \equiv X$, so the rule of logical equivalence gives $P(\varphi \mid X) = p$, showing that $Q \subseteq P$.

Now assume $Q \subseteq P$. Let $(Y, \varphi, p) \in Q_0$. Choose $(X, \varphi, p) \in Q$ such that $Y = T_0 + \psi_X$. Then $P(\varphi \mid X) = p$. Note that $\widehat{X} \equiv Y$, so that $X \equiv Y \cup S_X$. If $S_X = \emptyset$, then $X \equiv Y$, so by the rule of logical equivalence, $P(\varphi \mid Y) = p$. Assume $S_X \neq \emptyset$. Let $\theta \in S_X \subseteq \tau(\widehat{Q}; \widehat{X})$. Then $(\widehat{X}, \theta, 1) \in \widehat{Q}$. But $\widehat{Q} \subseteq Q \subseteq P$, so $P(\theta \mid \widehat{X}) = 1$. Since $\widehat{X} \equiv Y$, the rule of logical equivalence implies $P(\theta \mid Y) = 1$. Since θ was arbitrary, deductive extension gives $Y \cup S_X \in \text{ante } P$ and $P(\cdot \mid Y) = P(\cdot \mid Y, S_X)$. Since $X \equiv Y \cup S_X$, we get $P(\varphi \mid Y, S_X) = p$. Therefore, $P(\varphi \mid Y) = p$, showing that $Q_0 \subseteq P$. \square

3.5.2 Intersections of inductive sets

We now get to the heart of the matter, which is the intersection of subsets of \mathcal{F}^{IS} . The next result shows that the closure properties are all preserved under such intersections. It is straightforward to verify that strong connectivity is also preserved. Hence, as a corollary, we find that the intersection of pre-theories (with a common root) is a pre-theory. After establishing that result, we give the proof of Theorem 3.3.7.

$\langle \text{T:intersect-cl} \rangle$ **Theorem 3.5.3.** *Let $\mathcal{C} \subseteq \mathfrak{P}\mathcal{F}^{\text{IS}}$ be nonempty.*

- (i) *If each set in \mathcal{C} is admissible, then $\bigcap \mathcal{C}$ is admissible.*
- (ii) *If each set in \mathcal{C} is entire, then $\bigcap \mathcal{C}$ is entire.*
- (iii) *If each set in \mathcal{C} is semi-closed, then $\bigcap \mathcal{C}$ is semi-closed.*
- (iv) *If each set in \mathcal{C} is closed, then $\bigcap \mathcal{C}$ is closed.*

Proof. Let $P = \bigcap \mathcal{C}$. Assume each set in \mathcal{C} is admissible. Suppose $(X, \varphi, p) \in P$, $X' \equiv X$, and $\varphi' \equiv_X \varphi$. Let $P' \in \mathcal{C}$, so that $P'(\varphi \mid X) = p$. Since P' is admissible, $P'(\varphi' \mid X') = p$. Since P' was arbitrary, $(X', \varphi', p) \in P$. Now suppose $(X', \varphi', p') \in P$. Choose $P' \in \mathcal{C}$. Then $P'(\varphi' \mid X') = p' = p$. Thus, P is admissible.

Now assume each set in \mathcal{C} is entire. Let $X \in \text{ante } P$ and $X \vdash \varphi$. Choose $\varphi' \in \mathcal{F}$ and $p \in [0, 1]$ such that $(X, \varphi', p) \in P$. Let $P' \in \mathcal{C}$. Then $P'(\varphi' \mid X) = p$, so that $X \in \text{ante } P'$. Since P' is entire and $X \vdash \varphi$, the rule of logical implication gives $P'(\varphi \mid X) = 1$. Since P' was arbitrary, we have $P(\varphi \mid X) = 1$, showing that P satisfies the rule of logical implication. Similar proofs show that P satisfies rules (R3)–(R7), and therefore, P is entire.

Assume each set in \mathcal{C} is semi-closed. Suppose $\overline{P}(\varphi \mid X) = p$ for every completion \overline{P} of P . Let $P' \in \mathcal{C}$ and let \overline{P}' be a completion of P' . Since $P \subseteq P' \subseteq \overline{P}'$, it follows that \overline{P}' is also a completion of P . Thus, $\overline{P}'(\varphi \mid X) = p$. Since \overline{P}' was arbitrary and P' is semi-closed, we have $P'(\varphi \mid X) = p$. Since P' was arbitrary, this gives $P(\varphi \mid X) = p$, and P is semi-closed.

Finally, assume each set in \mathcal{C} is closed. Suppose $S \subseteq \mathcal{F}$ is nonempty and $P(\theta \mid X) = 1$ for all $\theta \in S$. That is, $\emptyset \neq S \subseteq \tau(P; X)$. Let $P' \in \mathcal{C}$. Since $P \subseteq P'$, we have $\tau(P; X) \subseteq \tau(P'; X)$, so that $\emptyset \neq S \subseteq \tau(P'; X)$. Since P' is closed, deductive extension implies $X \cup S \in \text{ante } P'$ and $P'(\cdot \mid X, S) = P'(\cdot \mid X)$. Since P' was arbitrary, we have

$$\begin{aligned} P(\varphi \mid X, S) = p &\text{ iff } P'(\varphi \mid X, S) = p \text{ for all } P' \in \mathcal{C} \\ &\text{ iff } P'(\varphi \mid X) = p \text{ for all } P' \in \mathcal{C} \\ &\text{ iff } P(\varphi \mid X) = p. \end{aligned}$$

Since S is a nonempty subset of $\tau(P; X)$, it follows that $\tau(P; X)$ is nonempty, which implies $X \in \text{ante } P$. Hence, by the above, $X \cup S \in \text{ante } P$ and $P(\cdot \mid X, S) = P(\cdot \mid X)$, showing that P is closed. \square

$\langle \mathbf{C}:\text{intersect-cl} \rangle$ **Corollary 3.5.4.** *Let \mathcal{C}^0 be a nonempty set of pre-theories with common root T_0 and define $P_0 = \bigcap \mathcal{C}^0$. Then P_0 is a pre-theory with root T_0 .*

Proof. By Proposition 3.5.3, the set P_0 is semi-closed. Let $X \in \text{ante } P_0$. Choose $\varphi \in \mathcal{F}$ and $p \in [0, 1]$ such that $P_0(\varphi \mid X) = p$. Let $P \in \mathcal{C}^0$. Then $P(\varphi \mid X) = p$, so that $X \in \text{ante } P$. Since P is strongly connected with root T_0 , we may choose $\psi \in \mathcal{F}$ such that $X \equiv T_0 + \psi$. Since X was arbitrary, P_0 is strongly connected with root T_0 . \square

Proof of Theorem 3.3.7. Let Q be consistent. Then Q is connected and we may choose an inductive theory P' such that $Q \subseteq P'$. Let T_0 be the root of Q . By Proposition 3.5.1, we have $T_0 \in \text{ante } P'$ and, if we define $P'_0 = P' \downarrow_{T_0}$, then P'_0 is a pre-theory with root T_0 . By Theorem 3.5.2, we may choose strongly connected Q_0 with root T_0 such that $Q_0 \equiv Q$. Let \mathcal{C}^0 be the set of all pre-theories with root T_0 that contain Q_0 .

Since $Q_0 \equiv Q$ and $Q \subseteq P'$, we have $Q_0 \subseteq P'$. Since Q_0 is strongly connected with root T_0 , this gives $Q_0 \subseteq P'_0$. Hence, \mathcal{C}^0 is nonempty. Let $P_0 = \bigcap \mathcal{C}^0$. Corollary 3.5.4 implies P_0 is a pre-theory with root T_0 . Define $P = \mathbf{P}(P_0)$. Since Q_0 is a subset of every element of \mathcal{C}^0 , it follows that $Q_0 \subseteq P_0 \subseteq P$. Therefore, since $Q_0 \equiv Q$, we have $Q \subseteq P$.

To show that P is the smallest such inductive theory, let P'' be an arbitrary inductive theory with $Q \subseteq P''$. As above, if $P''_0 = P'' \downarrow_{T_0}$, then P''_0 is a pre-theory with root T_0 and $Q_0 \subseteq P''_0$. Hence, $P''_0 \in \mathcal{C}^0$, so that $P_0 \subseteq P''_0$, which implies $P \subseteq \mathbf{P}(P''_0)$. But $P''_0 \subseteq P''$, so by Theorem 3.3.4, we have $\mathbf{P}(P''_0) \subseteq P''$, and therefore, $P \subseteq P''$. \square

3.5.3 A converse to the rule of logical implication

Having proved Theorem 3.3.7, we now have the notation \mathbf{P}_Q at our disposal. With this notation, we are able to give a more natural characterization of the equivalence of two subsets of \mathcal{F}^{IS} . We also give a definition that we will need later, and provide a partial converse to the rule of logical implication.

$\langle \mathbf{P}:\text{ind-equiv-char} \rangle$ **Proposition 3.5.5.** *Let $Q, Q' \subseteq \mathcal{F}^{\text{IS}}$ be connected. Then $Q \equiv Q'$ if and only if either both are inconsistent or both are consistent and $\mathbf{P}_Q = \mathbf{P}_{Q'}$.*

Proof. Assume $Q \equiv Q'$. Then $Q \subseteq P$ if and only if $Q' \subseteq P$ for all inductive theories P . Hence, Q is consistent if and only if Q' is consistent. Suppose both are consistent. Since $Q \subseteq \mathbf{P}_Q$, we have $Q' \subseteq \mathbf{P}_Q$, which implies $\mathbf{P}_{Q'} \subseteq \mathbf{P}_Q$. Similarly, $\mathbf{P}_Q \subseteq \mathbf{P}_{Q'}$, and therefore, $\mathbf{P}_Q = \mathbf{P}_{Q'}$.

For the converse, if both are inconsistent, then neither can be extended to an inductive theory, so they are vacuously equivalent. Assume both are consistent and $\mathbf{P}_Q = \mathbf{P}_{Q'}$. Let P be an inductive theory with $Q \subseteq P$. Then $\mathbf{P}_Q \subseteq P$. Thus, $Q' \subseteq \mathbf{P}_{Q'} = \mathbf{P}_Q \subseteq P$. Similarly, $Q' \subseteq P$ implies $Q \subseteq P$, showing that $Q \equiv Q'$. \square

If Q is consistent, we define $T(Q) = \tau(\mathbf{P}_Q)$. We also denote $T(Q)$ by T_Q . By Proposition 3.2.13, the set T_Q is a deductive theory. We call T_Q the *deductive theory determined by Q* .

Note that $\varphi \in T_Q$ if and only if $Q \vdash (X, \varphi, 1)$ for all $X \in \text{ante } \mathbf{P}_Q$. In other words, T_Q is the set of formulas which, under \mathbf{P}_Q , have probability one, regardless of the antecedent. Informally, T_Q represents the deductive hypotheses that are implicit in the set Q .

(T:log-impl-iff) **Theorem 3.5.6.** *Suppose P is an inductive theory and let $X \in \text{ante } P$. Then $P(\varphi | X) = 1$ if and only if $X, T_P \vdash \varphi$.*

Proof. Let P be an inductive theory and $X \in \text{ante } P$. Then $P(\theta | X) = 1$ for all $\theta \in T_P$. Hence, by the rule of deductive extension, $X \cup T_P \in \text{ante } P$ and $P(\cdot | X) = P(\cdot | X, T_P)$.

Suppose $X, T_P \vdash \varphi$. Since $X \cup T_P \in \text{ante } P$, the rule of logical implication gives $P(\varphi | X, T_P) = 1$. By the above, $P(\varphi | X) = 1$. For the converse, suppose $P(\varphi | X) = 1$. Let T_0 be the root of P , so that $P = \mathbf{P}(P_0)$, where $P_0 = P \downarrow_{T_0}$. We then have $X \equiv T^X + \psi_X$ and $P_0(\varphi | T_0, \psi_X) = 1$. By the rule of material implication, $P_0(\psi_X \rightarrow \varphi | T_0) = 1$, which implies $\psi_X \rightarrow \varphi \in T_P$. Therefore, $X, T_P \vdash \psi_X, \psi_X \rightarrow \varphi \vdash \varphi$. \square

3.5.4 Inductive conditions

At this point, we have proven all the results in Sections 3.3.4 and 3.3.5. We have therefore fully established and justified the notation $Q \vdash (X, \varphi, p)$. Informally, we think of $Q \vdash (X, \varphi, p)$ as representing a process of derivation, where we take the inductive statements in Q as our hypotheses, then apply the nine rules of inductive inference to derive (X, φ, p) . But every hypothesis in Q is a precise inductive statements of the form $Q(\eta | Y) = q$. We are often interested in using more general hypotheses, such as $Q(\eta | Y) > q$. Or we may wish to hypothesize that Q and everything it entails satisfies a certain symmetry condition. To allow for these more general hypotheses, we define the following.

An *inductive condition* is a collection \mathcal{C} of inductive theories with a common root. An inductive condition \mathcal{C} is said to be *consistent* if $\mathcal{C} \neq \emptyset$. If each $P \in \mathcal{C}$ has root T_0 , then we call T_0 the *root* of \mathcal{C} .

For instance, \mathcal{C} might be a collection of inductive theories P , each satisfying $P(\eta | Y) > q$. Or \mathcal{C} might be a collection where each member satisfies a given symmetry property. If we wish to simply assume Q , without any generalizations, we can also do that with an inductive condition. Namely, if Q is connected with root T_0 , then let $\mathcal{C}(Q)$, which we also denote by \mathcal{C}_Q , be defined as the set of inductive theories P with root T_0 such that $Q \subseteq P$. Then \mathcal{C}_Q also has root T_0 and \mathcal{C}_Q is consistent if and only if Q is consistent. Moreover, by Theorem 3.3.7, we have $\mathbf{P}_Q \in \mathcal{C}_Q$ and $\mathbf{P}_Q \subseteq P$ for all $P \in \mathcal{C}_Q$. Hence, $\mathbf{P}_Q = \bigcap \mathcal{C}_Q$. If we identify \mathcal{C}_Q with Q , then we can say that \mathcal{C}_Q generates the inductive theory $\bigcap \mathcal{C}_Q$.

Generalizing this to arbitrary inductive conditions is not straightforward. The problem, as mentioned at the beginning of this section, is that the

intersection of inductive theories is not necessarily an inductive theory. What we will show, however, is that if \mathcal{C} is an inductive condition with root T_0 , then there is a largest inductive theory with root T_0 contained in $\bigcap \mathcal{C}$. If $\mathcal{C} = \mathcal{C}_Q$, then that largest theory is \mathbf{P}_Q .

Let \mathcal{C} be an inductive condition with root T_0 . Define $\mathcal{C}^0 = \{P \downarrow_{T_0} \mid P \in \mathcal{C}\}$. By Corollary 3.5.4, if \mathcal{C} is consistent, then $\bigcap \mathcal{C}^0$ is a pre-theory with root T_0 . Hence, may define $\mathbf{P}(\mathcal{C}) = \mathbf{P}(\bigcap \mathcal{C}^0)$. We also denote $\mathbf{P}(\mathcal{C})$ by $\mathbf{P}_{\mathcal{C}}$, and call this the *inductive theory generated by \mathcal{C}* . The next result shows that the inductive theory generated by \mathcal{C} is indeed the largest inductive theory contained in $\bigcap \mathcal{C}$.

(T:theory-gen-IC-defn) **Theorem 3.5.7.** *Let \mathcal{C} be a consistent inductive condition. Then $\mathbf{P}_{\mathcal{C}} \subseteq \bigcap \mathcal{C}$. Moreover, if P is an inductive theory with the same root as \mathcal{C} such that $P \subseteq \bigcap \mathcal{C}$, then $P \subseteq \mathbf{P}_{\mathcal{C}}$.*

Proof. Let \mathcal{C} be a consistent inductive condition with root T_0 . Let $P' \in \mathcal{C}$. Then $\bigcap \mathcal{C}^0 \subseteq \bigcap \mathcal{C} \subseteq P'$, which implies $\mathbf{P}_{\mathcal{C}} \subseteq P'$. Since P' was arbitrary, this shows $\mathbf{P}_{\mathcal{C}} \subseteq \bigcap \mathcal{C}$.

Now suppose P is an inductive theory with root T_0 such that $P \subseteq \bigcap \mathcal{C}$. Then $P = \mathbf{P}(P_0)$, where $P_0 = P \downarrow_{T_0}$. Let $P'_0 \in \mathcal{C}^0$ be arbitrary, and choose $P' \in \mathcal{C}$ such that $P'_0 = P' \downarrow_{T_0}$. Since $P' \in \mathcal{C}$, it follows that $P \subseteq P'$, which implies $P_0 \subseteq P'_0$. Since P'_0 was arbitrary, this gives $P_0 \subseteq \bigcap \mathcal{C}^0$. Therefore, $P = \mathbf{P}(P_0) \subseteq \mathbf{P}(\bigcap \mathcal{C}^0) = \mathbf{P}_{\mathcal{C}}$. \square

As noted earlier, if $Q \subseteq \mathcal{F}^{\text{IS}}$ is consistent, then $\mathbf{P}_Q = \bigcap \mathcal{C}_Q$. Hence, by Theorem 3.5.7, we have $\mathbf{P}(Q) = \mathbf{P}(\mathcal{C}_Q)$.

If \mathcal{C} is an inductive condition and $\mathbf{P}_{\mathcal{C}} \in \mathcal{C}$, then we say the condition \mathcal{C} is *determinate*, otherwise \mathcal{C} is *indeterminate*. Note that if $Q \subseteq \mathcal{F}^{\text{IS}}$ is consistent, then \mathcal{C}_Q is determinate.

If \mathcal{C} is an inductive condition and $(X, \varphi, p) \in \mathcal{F}^{\text{IS}}$, we write $\mathcal{C} \vdash (X, \varphi, p)$ to mean that \mathcal{C} is consistent and $\mathbf{P}_{\mathcal{C}}(\varphi \mid X) = p$. Since $\mathbf{P}(Q) = \mathbf{P}(\mathcal{C}_Q)$, we have that $Q \vdash (X, \varphi, p)$ if and only if $\mathcal{C}_Q \vdash (X, \varphi, p)$, so that this new use of the turnstile symbol is an extension of our previous use.

If \mathcal{C} and \mathcal{C}' are inductive conditions with the same root, then $\mathcal{C}, \mathcal{C}' \vdash (X, \varphi, p)$ means $\mathcal{C} \cap \mathcal{C}' \vdash (X, \varphi, p)$. In particular, if we identify Q and \mathcal{C}_Q , then $Q, \mathcal{C} \vdash (X, \varphi, p)$ means $\mathcal{C}_Q, \mathcal{C} \vdash (X, \varphi, p)$. Lastly, we use $\mathcal{C}, X \vdash \varphi$ as shorthand for $\mathcal{C} \vdash (X, \varphi, 1)$. For example, $Q, \mathcal{C}, X \vdash \varphi$ means that Q and \mathcal{C} have the same root, $\mathcal{C}_Q \cap \mathcal{C}$ is consistent (that is, nonempty), and, if $P = \mathbf{P}(\mathcal{C}_Q \cap \mathcal{C})$, then $P(\varphi \mid X) = 1$.

Finally, if \mathcal{C} is a consistent inductive condition, we define $T(\mathcal{C}) = T(\mathbf{P}_{\mathcal{C}})$. We also denoted $T(\mathcal{C})$ by $T_{\mathcal{C}}$.

(P:ded-th-ind-cond) **Proposition 3.5.8.** *If \mathcal{C} is a consistent inductive condition, then $T_{\mathcal{C}} = \bigcap \{T_P \mid P \in \mathcal{C}\}$.*

Proof. Let \mathcal{C} be a consistent inductive condition with root T_0 . Let $P'_0 = \bigcap \mathcal{C}^0$ and $P' = \mathbf{P}_{\mathcal{C}}$, so that $P' = \mathbf{P}(P'_0)$. Suppose $\theta \in T_{\mathcal{C}} = T_{P'}$. Then $P'(\theta \mid T_0) = 1$, so by Theorem 3.5.7, we have $P(\theta \mid T_0) = 1$ for all $P \in \mathcal{C}$. By Proposition 3.4.5, we have $\theta \in T_P$ for all $P \in \mathcal{C}$, so that $\theta \in \bigcap \{T_P \mid P \in \mathcal{C}\}$.

Conversely, suppose $\theta \in \bigcap \{T_P \mid P \in \mathcal{C}\}$. Then $P_0(\theta \mid T_0) = 1$ for all $P_0 \in \mathcal{C}^0$. Hence, $P'_0(\theta \mid T_0) = 1$. As above, Proposition 3.4.5 gives $\theta \in T_{P'}$. \square

Recall Theorem 3.5.6, which gives a partial converse to the rule of logical implication. The following result reformulates that in terms of our new notation and shows that $\mathcal{C}, X \vdash \varphi$ can be rewritten in terms of the classical derivability relation from Section 3.1.

Proposition 3.5.9. *Suppose \mathcal{C} is a consistent inductive condition and let $X \in \text{ante } \mathbf{P}_{\mathcal{C}}$. Then $\mathcal{C}, X \vdash \varphi$ if and only if $T_{\mathcal{C}}, X \vdash \varphi$.*

Proof. Let \mathcal{C} be consistent and let $X \in \text{ante } \mathbf{P}_{\mathcal{C}}$. Note that $\mathcal{C}, X \vdash \varphi$ if and only if $\mathbf{P}_{\mathcal{C}}(\varphi \mid X) = 1$. Also note that $T_{\mathcal{C}} = T(\mathbf{P}_{\mathcal{C}})$. Hence, the result follows immediately from Theorem 3.5.6. \square

We conclude this section with a result that we will need in Chapter 4.

(P:chop-off-root) **Proposition 3.5.10.** *Let P be an inductive theory with root T_0 and let $T'_0 \in [T_0, T_P]$. Let $P'_0 = P \downarrow_{T'_0}$. Then P'_0 is a pre-theory with root T'_0 . Moreover, if $P' = \mathbf{P}(P'_0)$, then $T_{P'} = T_P$ and $P' = P \downarrow_{[T'_0, T_P]}$.*

Proof. By Theorem 3.4.7, since P'_0 is strongly connected, we have that P'_0 is a pre-theory with root T'_0 . Let $P' = \mathbf{P}(P'_0)$. We first show that $T_{P'} = T_P$. By Proposition 3.4.5, it suffices to show that $P'(\theta \mid T'_0) = 1$ if and only if $P(\theta \mid T_0) = 1$. Suppose $P'(\theta \mid T'_0) = 1$. Then $P(\theta \mid T'_0) = 1$. But $T'_0 \subseteq T_P$, so by the rule of deductive extension, we have $P(\theta \mid T_0) = 1$. Conversely, suppose $P(\theta \mid T_0) = 1$. Again by the rule of deductive extension, we obtain $P(\theta \mid T'_0) = 1$, and so $P'(\theta \mid T'_0) = 1$.

It remains to show that $P' = P \downarrow_{[T'_0, T_P]}$. Since $P'_0 \subseteq P$, Theorem 3.3.4 implies $P' \subseteq P$. Moreover, every $X \in \text{ante } P'$ satisfies $X \leftrightarrow [T'_0, T_{P'}] = [T'_0, T_P]$. Hence, $P' \subseteq P \downarrow_{[T'_0, T_P]}$.

Conversely, suppose $(X, \varphi, p) \in P \downarrow_{[T'_0, T_P]}$. Then $P(\varphi \mid X) = p$ and $X \leftrightarrow [T'_0, T_P]$. Write $T(X) = T + \psi$, where $T \in [T'_0, T_P] \subseteq [T_0, T_P]$. Then $p = P(\varphi \mid X) = P_0(\varphi \mid T_0, \psi)$. Since $P_0 \subseteq P$, we have $P(\varphi \mid T_0, \psi) = p$. But $T_0 \subseteq T'_0 \subseteq \tau(P)$, so by the rule of deductive extension, we have $P(\varphi \mid T'_0, \psi) = p$, and this implies $P'_0(\varphi \mid T'_0, \psi) = p$. Since $T_P = T_{P'}$, we have $T \in [T'_0, T_{P'}]$. Therefore $P'(\varphi \mid X) = p$, and this shows $P \downarrow_{[T'_0, T_P]} \subseteq P'$. \square

Chapter 4

Propositional Models

`<Ch:prop-models>` The logical relationships between deductive and inductive statements in \mathcal{F} and \mathcal{F}^{IS} are described by the derivability relation \vdash developed in Chapter 3. This relation is a kind of calculus, based on a set of inferential rules. The rules themselves depend only on the syntax of the statements. No interpretation or meaning is given to them, and no such meaning is necessary to describe these logical relationships.

We can, however, use meanings and interpretations to investigate the logical relationships between statements. When we do this, we arrive at a different relation, called the consequence relation, and denoted by \models . When we study logical relationships using \vdash , we are studying the syntactics of the logic. When we use \models , we are studying the semantics.

Meanings are assigned to statements using models. The classical type of propositional model (which we call a strict model) is one that simply assigns a truth value (0 or 1) to each propositional variable. Truth values then propagate to every sentence in \mathcal{F} via the usual interpretations of \neg and \wedge . These strict models originated with Wittgenstein (see [33, Satz 4.31]) and can be visualized as rows in what we now call a truth table.

In Wittgenstein's view, a strict model represents a logically possible state of the world. Since we are interested in modeling inductive inference, we wish to model degrees of uncertainty about the state of the world. For us, then, a model will be a collection of strict models, together with a set of weights whose relative magnitudes represent relative degrees of uncertainty. Without loss of generality, we may assume these weights add up to one. In other words, we will define a model to be a probability measure on a set of strict models.

Models are used to define the consequence relation \models . On the deductive side, we say that $X \models \varphi$ if, in every model that satisfies X , the sentence φ is also satisfied. A completely analogous definition holds on the inductive side when we write $P \models (X, \varphi, p)$.

In Section 4.1, we define models and we define the consequence relation for deductive statements. That is, we define what it means to say that $X \models \varphi$, when $X \subseteq \mathcal{F}$ and $\varphi \in \mathcal{F}$. We prove that, together, \vdash and \models form a sound

logical system, meaning that $X \vdash \varphi$ implies $X \models \varphi$. In other words, if we can use X to prove φ , then φ is satisfied in every model of X . We also prove that this logical system is complete, meaning that $X \models \varphi$ implies $X \vdash \varphi$. In other words, if φ is satisfied in every model of X , then there is a proof of φ from X . Together, soundness and completeness show that $X \vdash \varphi$ if and only if $X \models \varphi$. In particular, by Theorem 3.1.10, the consequence relation is σ -compact.

These results should be contrasted with the usual approach to (deductive) semantics in \mathcal{F} . One can define a strict consequence relation using strict models. It is well known (see Examples 4.4.5 and 4.4.6) that both completeness and σ -compactness fail in that case. By overlaying our semantics with a probability measure, we recover both properties.

In Section 4.2, we extend the consequence relation to inductive statements, and prove both the soundness and completeness of this extension. These results finalize our description of probability as logic. Probability is a system of inference on inductive statements. It contains classical propositional logic as a special case and extends it in two directions: from deductive to inductive, and from finite conjunctions to countably infinite conjunctions. It has both semantics and a syntactic calculus. (In Chapter 5, we will repeat these constructions in a predicate language, showing that probability, as inductive logic, also extends first-order logic in both these directions.)

In Section 4.3.1, we address the relationship between this logical system of probability and modern, measure-theoretic probability theory. Modern probability has its origin in Kolmogorov's 1933 manuscript, *Foundations of the Theory of Probability* [20]. Therein, Kolmogorov lays out what he calls the axioms of probability. Today, those axioms take the form of a definition, namely, the definition of a probability space: a measure space with total mass one. The foundation of modern probability, therefore, is the probability space.

For us, the probability space is the foundation of our semantics. In Theorem 4.3.1, we show that every probability space is isomorphic to a semantic model in our logical system. The proof of Theorem 4.3.1 exhibits a natural mapping from the outcomes and events of the given probability space to strict models and sentences, respectively. With this result, we see that all of modern, measure-theoretic probability is embedded in our logical system. Probability theory as we know it today is simply the semantics of a larger system of logical reasoning. (In Chapter 5, we will extend this embedding to include random variables. See, for instance, Section 5.4.1 and the table of correspondences immediately preceding Section 5.4.2.)

Generally speaking, there is a difference between saying that a model satisfies a set of statements, and saying that it characterizes them. To say that it characterizes them is to say that they are the only statements satisfied by the model. In that case, the model is a perfect semantic reflection of the given set of statements. Finding a model that characterizes a set of statements allows us to see the logical structure of that set in a single semantic model.

In Theorem 4.2.4, we show that models (that is, probability spaces), are only able to characterize complete inductive theories. The structure of an incomplete inductive theory cannot be represented by a probability space. Proposition

4.1.16 tells us why. Namely, if we start with a collection of inductive statements and conditions, and then draw all available inferences, we are not led to a σ -algebra, but rather to a Dynkin system. Consequently, we see that Dynkin systems arise naturally and organically in the study of inductive inference. This fact may offer some insight into why Dynkin's π - λ theorem—a purely measure-theoretic result—features so much more prominently in probability theory than it does in analysis.

In the rest of Section 4.3 and in Section 4.4, we present several examples, illustrating and applying many of the ideas presented in both Chapters 3 and 4.

Finally, in Section 4.5, we introduce the idea of (inductive) independence—a purely logical and syntactic notion—and then show that it is semantically characterized by the usual product formula from measure theory. We then present two examples to illustrate its use.

4.1 Models and deductive semantics

(S:models_ded_sem)

4.1.1 Truth assignments

Recall that \mathbf{B} denotes the Boolean σ -algebra $\{0, 1\}$, whose partial order is the usual \leq . The elements 0 and 1 are called *truth values*. If S is a set, then a function $\nu : S \rightarrow \mathbf{B}$ is an assignment of truth values to the elements of S . The set of all such functions is denoted by \mathbf{B}^S .

Given an element $s \in S$, we define the *projection* $\pi_s : \mathbf{B}^S \rightarrow \mathbf{B}$ by $\pi_s \nu = \nu s$. Let $\mathcal{B} = \mathfrak{P} \mathbf{B}$. Then \mathcal{B}^S denotes the product σ -algebra. That is, \mathcal{B}^S is the smallest σ -algebra on \mathbf{B}^S such that each π_s is $(\mathcal{B}^S, \mathcal{B})$ -measurable. In symbols,

$$\mathcal{B}^S = \sigma(\{\pi_s \mid s \in S\}) = \sigma(\{\pi_s^{-1} A \mid s \in S, A \in \mathcal{B}\}).$$

A subset of \mathbf{B}^S is called a *cylinder set* if it has the form

$$\pi_{s_1}^{-1} A_1 \cap \cdots \cap \pi_{s_n}^{-1} A_n$$

for some $s_1, \dots, s_n \in S$ and $A_1, \dots, A_n \in \mathcal{B}$. Equivalent, a cylinder set is a set of the form

$$\{\nu \in \mathbf{B}^S \mid \nu s_1 = x_1, \dots, \nu s_n = x_n\}.$$

for some $s_1, \dots, s_n \in S$ and $x_1, \dots, x_n \in \mathbf{B}$. The σ -algebra \mathcal{B}^S is also generated by the collection of cylinder sets. Note that the collection of cylinder sets is a π -system. That is, it is closed under intersections.

If we say that a function $f : \mathbf{B}^S \rightarrow \mathbf{B}$ is measurable, we mean that it is $(\mathcal{B}^S, \mathcal{B})$ -measurable. Note that $\neg f = 1 - f$ and $\bigwedge_n f_n = \inf_n f_n$. Hence, if f and f_n are all measurable, then so are $\neg f$ and $\bigwedge_n f_n$. Also note that every function $f : \mathbf{B}^S \rightarrow \mathbf{B}$ has the form $f = 1_B$ for some $B \subseteq \mathbf{B}^S$, and f is measurable if and only if $B \in \mathcal{B}^S$.

(L:ctble-ary)

Lemma 4.1.1. *Let $R \subseteq S$, let $h : \mathbf{B}^R \rightarrow \mathbf{B}$, and define $f : \mathbf{B}^S \rightarrow \mathbf{B}$ by $f \nu = h(\nu|_R)$. If h is measurable, then f is measurable.*

Proof. Let $R \subseteq S$ and define

$$\Sigma = \{B \in \mathcal{B}^R \mid \nu \mapsto 1_B(\nu|_R) \text{ is } (\mathcal{B}^S, \mathcal{B})\text{-measurable}\}.$$

It suffices to show that $\Sigma = \mathcal{B}^R$. Since constant functions are measurable, we have $\emptyset \in \Sigma$. Since $1_{B^c} = \neg 1_B$, it follows that Σ is closed under complements. Let $\{B_n\}_{n=1}^\infty \subseteq \Sigma$ and define $B = \bigcap_n B_n$. Then $1_B = \bigwedge_n 1_{B_n}$, so that $B \in \Sigma$, and Σ is closed under countable intersection. Therefore, Σ is a σ -algebra.

Now let

$$B = \{\nu \in \mathcal{B}^R \mid \nu s_1 = x_1, \dots, \nu s_n = x_n\}$$

be a cylinder set in \mathcal{B}^R . Define $h : \mathcal{B}^S \rightarrow \mathcal{B}$ by $h\nu = 1_B(\nu|_R)$. Then

$$h^{-1}1 = \{\nu \in \mathcal{B}^S \mid \nu s_1 = x_1, \dots, \nu s_n = x_n\}$$

is a cylinder set in \mathcal{B}^S . Hence, Σ contains the cylinder sets in \mathcal{B}^R . Since \mathcal{B}^R is the smallest σ -algebra containing the cylinder sets, we have $\Sigma = \mathcal{B}^R$. \square

Let $R \subseteq S$. A measurable function $f : \mathcal{B}^S \rightarrow \mathcal{B}$ is said to be R -ary if there exists a measurable $h : \mathcal{B}^R \rightarrow \mathcal{B}$ such that $f\nu = h(\nu|_R)$ for all $\nu \in \mathcal{B}^S$.

$\langle \text{C:ctble-ary} \rangle$ **Corollary 4.1.2.** *Let $R \subseteq U \subseteq S$. If $f : \mathcal{B}^S \rightarrow \mathcal{B}$ is R -ary, then f is U -ary.*

Proof. Let $f : \mathcal{B}^S \rightarrow \mathcal{B}$ be R -ary. Choose measurable $h : \mathcal{B}^R \rightarrow \mathcal{B}$ such that $f\nu = h(\nu|_R)$ for all $\nu \in \mathcal{B}^S$. Define $g : \mathcal{B}^U \rightarrow \mathcal{B}$ by $g\nu = h(\nu|_R)$. Lemma 4.1.1 implies g is measurable. Moreover, for any $\nu \in \mathcal{B}^S$, we have $g(\nu|_U) = h((\nu|_U)|_R) = h(\nu|_R) = f\nu$. \square

$\langle \text{P:ctble-ary} \rangle$ **Proposition 4.1.3.** *Let S be a set. Then every measurable $f : \mathcal{B}^S \rightarrow \mathcal{B}$ is R -ary for some countable $R \subseteq S$.*

Proof. Let

$$\Sigma = \{B \in \mathcal{B}^S \mid 1_B \text{ is } R\text{-ary for some countable } R \subseteq S\}.$$

It suffices to show that $\Sigma = \mathcal{B}^S$. Clearly, 1_\emptyset is \emptyset -ary, so that $\emptyset \in \Sigma$. Since $1_{B^c} = \neg 1_B$, it follows that Σ is closed under complements. Let $\{B_n\}_{n=1}^\infty \subseteq \Sigma$ and define $B = \bigcap_n B_n$. For each n , choose countable $R_n \subseteq S$ such that 1_{B_n} is R_n -ary. Let $R = \bigcup_n R_n$. Then R is countable and, by Corollary 4.1.2, it follows that 1_{B_n} is R -ary for all n . Choose measurable $h_n : \mathcal{B}^R \rightarrow \mathcal{B}$ such that $1_{B_n}\nu = h_n(\nu|_R)$ for all $\nu \in \mathcal{B}^S$, and define $h = \bigwedge_n h_n$. Then h is measurable and

$$1_B\nu = \bigwedge_n 1_{B_n}\nu = \bigwedge_n h_n(\nu|_R) = h(\nu|_R),$$

so that $B \in \Sigma$. Hence, Σ is closed under countable intersections, and Σ is therefore a σ -algebra.

Now let

$$B = \{\nu \in \mathcal{B}^S \mid \nu s_1 = x_1, \dots, \nu s_n = x_n\}$$

be a cylinder set in \mathcal{B}^S . Then B is R -ary, where $R = \{s_1, \dots, s_n\}$, showing that Σ contains the cylinder sets. Since \mathcal{B}^S is the smallest σ -algebra containing the cylinder sets, we have $\Sigma = \mathcal{B}^S$. \square

4.1.2 Strict models and Boolean functions

A *strict model* is a function $\omega : PV \rightarrow \mathbf{B}$. That is, a strict model is an assignment of truth values to each propositional formula. The set of all strict models is \mathbf{B}^{PV} . The domain of a strict model can be uniquely extended to all of \mathcal{F} by formula recursion. That is, $\omega \neg \varphi = \neg \omega \varphi$ and $\omega \bigwedge \Phi = \bigwedge_{\varphi \in \Phi} \omega \varphi$. We call $\omega \varphi$ the *truth value* of φ in the strict model ω .

The definition of a strict model depends on the choice of propositional variables PV , which in turn determine the language \mathcal{F} . When we wish to emphasize this fact, we will call ω a *strict model in \mathcal{F}* .

A *Boolean function* is a measurable function $f : \mathbf{B}^{PV} \rightarrow \mathbf{B}$. We say that a formula φ *represents* a Boolean function f if $\omega \varphi = f\omega$ for all strict models ω .

(P:Boolean-funcs) **Proposition 4.1.4.** *Every formula represents a unique Boolean function. Conversely, every Boolean function is represented by a formula.*

Proof. Let $\varphi \in \mathcal{F}$ and define $f_\varphi : \mathbf{B}^{PV} \rightarrow \mathbf{B}$ by $f_\varphi \omega = \omega \varphi$. To show that f_φ is a Boolean function, we must show that it is measurable. This follows by formula induction since $f_{\mathbf{r}} = \pi_{\mathbf{r}}$, $f_{\neg \varphi} = \neg f_\varphi$, and $f_{\bigwedge \Phi} = \bigwedge_{\varphi \in \Phi} f_\varphi$.

Now let

$$\Sigma = \{B \in \mathcal{B}^{PV} \mid 1_B \text{ is represented by a formula}\}.$$

It suffices to show that $\Sigma = \mathcal{B}^{PV}$. If we fix $\mathbf{r} \in PV$, then 1_\emptyset is represented by the formula $\mathbf{r} \wedge \neg \mathbf{r}$, so $\emptyset \in \Sigma$. If φ represents 1_B , then $\neg \varphi$ represents $\neg 1_B = 1_{B^c}$, so Σ is closed under complements. And if φ_n represents 1_{B_n} , then $\bigwedge_n \varphi_n$ represents $\bigwedge_n 1_{B_n} = 1_{\bigcap_n B_n}$, so that Σ is closed under countable intersections, and Σ is therefore a σ -algebra.

Now let

$$B = \{\omega \in \mathbf{B}^{PV} \mid \omega \mathbf{r}_1 = x_1, \dots, \omega \mathbf{r}_n = x_n\}$$

be a cylinder set. Recall the notation $\varphi^1 = \varphi$ and $\varphi^0 = \neg \varphi$, and note that $\omega \varphi = x$ if and only if $\omega \varphi^x = 1$. Hence, 1_B is represented by the formula $\mathbf{r}_1^{x_1} \wedge \dots \wedge \mathbf{r}_n^{x_n}$, so that Σ contains the cylinder sets, and therefore $\Sigma = \mathcal{B}^{PV}$. \square

By Proposition 4.1.3, every Boolean function is Π -ary for some countable set of propositional variables $\Pi \subseteq PV$.

(P:Boolean-func-Pi-ary) **Proposition 4.1.5.** *Let $\varphi \in \mathcal{F}$ and let f be the Boolean function that it represents. Then f is Π -ary, where $\Pi = PV \cap \text{Sf } \varphi$ is the countable set of propositional variables that appear in φ .*

Proof. Let $\varphi \in \mathcal{F}$ represent f_φ and let $\Pi_\varphi = PV \cap \text{Sf } \varphi$. We will show that f_φ is Π_φ -ary by induction on φ .

If $\mathbf{r} \in PV$, then $\Pi_{\mathbf{r}} = \{\mathbf{r}\}$ and $f_{\mathbf{r}} = \pi_{\mathbf{r}}$ is $\{\mathbf{r}\}$ -ary. Suppose f_φ is Π_φ -ary. Then $f_{\neg \varphi} = \neg f_\varphi$ is also Π_φ -ary. Since $\text{Sf } \varphi \subseteq \text{Sf } \neg \varphi$, we have $\Pi_\varphi \subseteq \Pi_{\neg \varphi}$. Hence, by Corollary 4.1.2, it follows that $f_{\neg \varphi}$ is $\Pi_{\neg \varphi}$ -ary.

Now let $\Phi \subseteq \mathcal{F}$ be countable and suppose f_θ is Π_θ -ary for all $\theta \in \Phi$. Define $\varphi = \bigwedge \Phi$. Note that $\Pi_\theta \subseteq \Pi_\varphi$ for each $\theta \in \Phi$. Hence, Corollary 4.1.2 implies that f_θ is Π_φ -ary for each $\theta \in \Phi$, and therefore $f_\varphi = \bigwedge_{\theta \in \Phi} f_\theta$ is also Π_φ -ary. \square

4.1.3 Models and satisfiability

An *inductive model*, or simply a *model*, is a probability space, $\mathcal{P} = (\Omega, \Sigma, \mathbb{P})$, where Ω is a set of strict models.

As with strict models, the definition of a model depends on the choice of propositional variables PV , which in turn determine the language \mathcal{F} . When we wish to emphasize this fact, we will call \mathcal{P} a *model in \mathcal{F}* .

If $X \subseteq \mathcal{F}$ and ω is a strict model, we say that ω *strictly satisfies* X , written $\omega \models X$, if $\omega\varphi = 1$ for all $\varphi \in X$. We write $\omega \models \varphi$ for $\omega \models \{\varphi\}$. A set $X \subseteq \mathcal{F}$ is *strictly satisfiable* if there is a strict model ω such that $\omega \models X$.

Let Ω be a set of strict models. For $\varphi \in \mathcal{F}$, let

$$\varphi_\Omega = \{\omega \in \Omega \mid \omega \models \varphi\}$$

be the set of strict models in Ω that strictly satisfy φ . More generally, for $X \subseteq \mathcal{F}$, we define $X_\Omega = \{\varphi_\Omega \mid \varphi \in X\}$.

The mapping $\varphi \mapsto \varphi_\Omega$ satisfies $(\neg\varphi)_\Omega = \varphi_\Omega^c$ and $(\bigwedge \Phi)_\Omega = \bigcap_{\varphi \in \Phi} \varphi_\Omega$. Similar relations hold for the shorthand operators. For instance $(\bigvee \Phi)_\Omega = \bigcup_{\varphi \in \Phi} \varphi_\Omega$, $(\varphi \rightarrow \psi)_\Omega = \varphi_\Omega^c \cup \psi_\Omega$, and $(\varphi \leftrightarrow \psi)_\Omega = (\varphi_\Omega \Delta \psi_\Omega)^c$.

Note that if $\varphi \in \mathcal{F}$ represents the Boolean function f , then $\varphi_\Omega = f^{-1}1$. Hence, by Proposition 4.1.4, if $\Omega = \mathbf{B}^{PV}$, then $\mathcal{B}^{PV} = \{\varphi_\Omega \mid \varphi \in \mathcal{F}\}$.

Let \mathcal{P} be a model and let $\varphi \in \mathcal{F}$. We say that \mathcal{P} *satisfies* φ , written $\mathcal{P} \models \varphi$, if $\overline{\mathbb{P}}\varphi_\Omega = 1$, where $(\Omega, \overline{\Sigma}, \overline{\mathbb{P}})$ is the completion of $(\Omega, \Sigma, \mathbb{P})$. For $X \subseteq \mathcal{F}$, we write $\mathcal{P} \models X$ to mean $\mathcal{P} \models \varphi$ for all $\varphi \in X$. Note that $\mathcal{P} \models \emptyset$ for every model \mathcal{P} . A set $X \subseteq \mathcal{F}$ is *satisfiable* if there is a model \mathcal{P} such that $\mathcal{P} \models X$. Note that if $X \subseteq X'$ and $\mathcal{P} \models X'$, then $\mathcal{P} \models X$.

(P:sig-pre-cpct) **Proposition 4.1.6.** *Let $X \subseteq \mathcal{F}$.*

(i) *If X is strictly satisfiable, then X is satisfiable.*

(ii) *If X is satisfiable and countable, then X is strictly satisfiable.*

Proof. Note that $\omega \models X$ if and only if $\mathcal{P} = (\{\omega\}, \{\emptyset, \{\omega\}\}, \delta_\omega) \models X$, which yields (i). For (ii), suppose X is satisfiable and countable. Let $\mathcal{P} = (\Omega, \Sigma, \mathbb{P})$ be a model that satisfies X . Then $\overline{\mathbb{P}}\bigcap_{\varphi \in X} \varphi_\Omega = 1$, so we may choose $\omega \in \bigcap_{\varphi \in X} \varphi_\Omega$, and this ω strictly satisfies X . \square

If \mathcal{P} is a model, we define

$$Th \mathcal{P} = \{\varphi \in \mathcal{F} \mid \mathcal{P} \models \varphi\}. \quad (4.1.1) \quad \boxed{\text{Th-sP-def}}$$

As we will see in Proposition 4.1.12, the set of formulas $Th \mathcal{P}$ is a consistent deductive theory. Note that if $\overline{\mathcal{P}} = (\Omega, \overline{\Sigma}, \overline{\mathbb{P}})$ is the completion of \mathcal{P} , then $Th \overline{\mathcal{P}} = Th \mathcal{P}$.

Let $\mathcal{P} = (\Omega, \Sigma, \mathbb{P})$ be a model with completion $(\Omega, \overline{\Sigma}, \overline{\mathbb{P}})$. Let $\overline{\Sigma}_{\mathcal{F}} = \overline{\Sigma} \cap \mathcal{B}^{PV}$ and let $\overline{\mathbb{P}}_{\mathcal{F}}$ be $\overline{\mathbb{P}}$ restricted to $\overline{\Sigma}_{\mathcal{F}}$. Then $\overline{\Sigma}_{\mathcal{F}}$ is a sub- σ -algebra of $\overline{\Sigma}$, so that $\mathcal{P}_{\mathcal{F}} = (\Omega, \overline{\Sigma}_{\mathcal{F}}, \overline{\mathbb{P}}_{\mathcal{F}})$ is also a model. For any $\varphi \in \mathcal{F}$, we have $\varphi_\Omega \in \overline{\Sigma}$ if and only if $\varphi_\Omega \in \overline{\Sigma}_{\mathcal{F}}$. Hence, all of the logical information in \mathcal{P} is contained in $\mathcal{P}_{\mathcal{F}}$. This is made precise in Proposition 4.1.8 below. We say that two models \mathcal{P} and \mathcal{Q} are *isomorphic (as models)*, denoted by $\mathcal{P} \simeq \mathcal{Q}$, if $\mathcal{P}_{\mathcal{F}} = \mathcal{Q}_{\mathcal{F}}$.

$\langle \text{L:prop-iso-thm} \rangle$ **Lemma 4.1.7.** *If \mathcal{P} is a model, then $(\mathcal{P}_{\mathcal{F}})_{\mathcal{F}} = \mathcal{P}_{\mathcal{F}}$.*

Proof. Let $\mathcal{P} = (\Omega, \Sigma, \mathbb{P})$ be a model. For notational simplicity, let $\Gamma = \overline{\Sigma}_{\mathcal{F}}$ and $\mathbb{Q} = \overline{\mathbb{P}}|_{\Gamma}$, so that $\mathcal{P}_{\mathcal{F}} = (\Omega, \Gamma, \mathbb{Q})$. We must show that $\overline{\Gamma}_{\mathcal{F}} = \Gamma$ and $\overline{\mathbb{Q}}_{\mathcal{F}} = \mathbb{Q}$. Since $\Gamma \subseteq \mathcal{B}^{PV}$, we have

$$\Gamma = \Gamma \cap \mathcal{B}^{PV} \subseteq \overline{\Gamma} \cap \mathcal{B}^{PV} = \overline{\Gamma}_{\mathcal{F}}.$$

Conversely, let $A \in \overline{\Gamma}_{\mathcal{F}} = \overline{\Gamma} \cap \mathcal{B}^{PV}$. Since $A \in \overline{\Gamma}$, we may write $A = B \cup N$, where $B \in \Gamma$, $N \subseteq F \in \Gamma$, and $\mathbb{Q}F = 0$. By the definition of \mathbb{Q} , this implies $\overline{\mathbb{P}}F = 0$. Now $\Gamma = \overline{\Sigma} \cap \mathcal{B}^{PV}$. Hence, $B \in \overline{\Sigma}$ and $F \in \overline{\Sigma}$. Since $N \subseteq F$ and $\overline{\mathbb{P}}F = 0$, we have $N \in \overline{\Sigma}$. Therefore, $A = B \cup N \in \overline{\Sigma}$. But $A \in \mathcal{B}^{PV}$ also. Hence, $A \in \overline{\Sigma} \cap \mathcal{B}^{PV} = \overline{\Sigma}_{\mathcal{F}} = \Gamma$, and this shows $\overline{\Gamma}_{\mathcal{F}} = \Gamma$.

By definition, $\overline{\mathbb{Q}}_{\mathcal{F}}$ is $\overline{\mathbb{Q}}$ restricted to $\overline{\Gamma}_{\mathcal{F}}$. But $\overline{\Gamma}_{\mathcal{F}} = \Gamma$, so we have that $\overline{\mathbb{Q}}_{\mathcal{F}} = \overline{\mathbb{Q}}|_{\Gamma} = \mathbb{Q}$. \square

$\langle \text{P:prop-iso-thm} \rangle$ **Proposition 4.1.8.** *For any model \mathcal{P} and any $\varphi \in \mathcal{F}$, we have $\mathcal{P} \models \varphi$ if and only if $\mathcal{P}_{\mathcal{F}} \models \varphi$.*

Proof. Suppose $\mathcal{P} \models \varphi$. Then $\varphi_{\Omega} \in \overline{\Sigma}$ and $\overline{\mathbb{P}}\varphi_{\Omega} = 1$. Since $\varphi_{\Omega} \in \mathcal{B}^{PV}$, we have $\varphi_{\Omega} \in \overline{\Sigma}_{\mathcal{F}}$ and $\overline{\mathbb{P}}_{\mathcal{F}}\varphi_{\Omega} = 1$. For the converse, let $\Gamma = \overline{\Sigma}_{\mathcal{F}}$ and $\mathbb{Q} = \overline{\mathbb{P}}|_{\Gamma}$, so that $\mathcal{P}_{\mathcal{F}} = (\Omega, \Gamma, \mathbb{Q})$. Suppose $\mathcal{P}_{\mathcal{F}} \models \varphi$. Then $\varphi_{\Omega} \in \overline{\Gamma}$ and $\overline{\mathbb{Q}}\varphi_{\Omega} = 1$. By Lemma 4.1.7, we have $\varphi_{\Omega} \in \Gamma$ and $\mathbb{Q}\varphi_{\Omega} = 1$. Since $\Gamma = \overline{\Sigma} \cap \mathcal{B}^{PV}$ and $\mathbb{Q} = \overline{\mathbb{P}}|_{\Gamma}$, this gives $\varphi_{\Omega} \in \overline{\Sigma}$ and $\overline{\mathbb{P}}\varphi_{\Omega} = 1$, so that $\mathcal{P} \models \varphi$. \square

Remark 4.1.9. If \mathcal{P} and \mathcal{Q} are isomorphic models, then $\mathcal{P} \models \varphi$ if and only if $\mathcal{Q} \models \varphi$ for all $\varphi \in \mathcal{F}$. This follows immediately from the definition of isomorphic models and Proposition 4.1.8.

4.1.4 Deductive consequence and soundness

We say $\varphi \in \mathcal{F}$ is a *consequence* of $X \subseteq \mathcal{F}$, or that X *entails* φ , which we denote by $X \models \varphi$, if, for all models \mathcal{P} such that $\mathcal{P} \models X$, we have $\mathcal{P} \models \varphi$. Note that if X is not satisfiable, then it is vacuously true that $X \models \varphi$ for all $\varphi \in \mathcal{F}$.

We write $\psi \models \varphi$ for $\{\psi\} \models \varphi$ and $\models \varphi$ for $\emptyset \models \varphi$. Note that $\models \varphi$ if and only if $\mathcal{P} \models \varphi$ for all models \mathcal{P} , which holds if and only if $\omega \models \varphi$ for all strict models ω . (If $\omega \models \varphi$ for all ω , then $\varphi_{\Omega} = \Omega$ in every model; conversely, if $\omega \not\models \varphi$, then $\mathcal{P} = (\{\omega\}, \{\emptyset, \{\omega\}\}, \delta_{\omega}) \not\models \varphi$.) We also write $X \models Y$ to mean that $X \models \varphi$ for all $\varphi \in Y$. Note here that $X \models Y$ if and only if $\mathcal{P} \models X$ implies $\mathcal{P} \models Y$ for all models \mathcal{P} .

A logical system is sound if every formula that is derivable from X is a consequence of X . The following theorem shows that our notion of deductive satisfiability yields a sound logical system, at least insofar as deductive inference is concerned.

$\langle \text{T:soundness} \rangle$ **Theorem 4.1.10 (Deductive soundness).** *Let $X \subseteq \mathcal{F}$ and $\varphi \in \mathcal{F}$. If $X \vdash \varphi$, then $X \models \varphi$.*

Proof. It suffices to show that (i)–(vi) in Definition 3.1.3 still hold when \vdash is replaced by \models . Conditions (i) and (ii) are trivial.

Suppose $X \models \bigwedge \Phi$. Let $\mathcal{P} = (\Omega, \Sigma, \mathbb{P}) \models X$. Then $\mathcal{P} \models \bigwedge \Phi$, which implies

$$(\bigwedge \Phi)_\Omega = \bigcap_{\theta \in \Phi} \theta_\Omega \in \bar{\Sigma}$$

and $\bar{\mathbb{P}} \bigcap_{\theta \in \Phi} \theta_\Omega = 1$. Thus, $\bar{\mathbb{P}} \bigcup_{\theta \in \Phi} \theta_\Omega^c = 0$. For each $\theta \in \Phi$, we have that θ_Ω^c is a subset of a null set. Hence, $\theta_\Omega^c \in \bar{\Sigma}$ and $\bar{\mathbb{P}} \theta_\Omega^c = 0$, implying $\theta_\Omega \in \bar{\Sigma}$ and $\bar{\mathbb{P}} \theta_\Omega = 1$. Therefore, $\mathcal{P} \models \theta$, showing that $X \models \theta$ and proving (iii). The proof of (iv) is similar.

For (v), suppose $X \models \varphi$ and $X \models \neg\varphi$, and assume there exists a model \mathcal{P} such that $\mathcal{P} \models X$. Then $\bar{\mathbb{P}} \varphi_\Omega = 1$ and $\bar{\mathbb{P}} \varphi_\Omega^c = 1$, a contradiction. Hence, X is not satisfiable, and so it is vacuously true that $X \models \psi$.

For (vi), suppose $X, \varphi \models \psi$, $X, \neg\varphi \models \psi$, and $X \not\models \psi$. Choose a model $\mathcal{P} = (\Omega, \Sigma, \mathbb{P})$ such that $\mathcal{P} \models X$ and $\mathcal{P} \not\models \psi$. If $\psi_\Omega \in \bar{\Sigma}$, then $\bar{\mathbb{P}} \psi_\Omega < 1$. Suppose $\psi_\Omega \notin \bar{\Sigma}$. Then $\mathbb{P}_* \psi_\Omega < \mathbb{P}^* \psi_\Omega \leq 1$. In this case, there exists a measure \mathbb{P}' on $(\Omega, \sigma(\Sigma \cup \{\psi_\Omega\}))$ such that $\mathbb{P}'|_\Sigma = \mathbb{P}$ and $\mathbb{P}' \psi_\Omega = \mathbb{P}_* \psi_\Omega$. In either case, $(\Omega, \Sigma, \mathbb{P})$ can be extended to a complete model $\mathcal{P}' = (\Omega, \Sigma', \mathbb{P}')$ in which $\psi_\Omega \in \Sigma'$ and $\bar{\mathbb{P}} \psi_\Omega < 1$. Therefore, $\mathcal{P}' \models X$ and $\mathcal{P}' \not\models \psi$.

By extending the model even further, we may assume $\varphi_\Omega \in \Sigma'$. Suppose $\mathbb{P}' \varphi_\Omega = 0$. Then $\mathcal{P}' \models \neg\varphi$, so by supposition, we have $\mathcal{P}' \models \psi$, a contradiction. Hence, $\mathbb{P}' \varphi_\Omega > 0$, and we may define a probability measure \mathbb{Q} on (Ω, Σ') by $\mathbb{Q} = \mathbb{P}'(\cdot | \varphi_\Omega)$, and then define the model $\mathcal{Q} = (\Omega, \Sigma', \mathbb{Q})$.

Since $\mathbb{Q} \varphi_\Omega = 1$, we have $\mathcal{Q} \models \varphi$. Also, if $A \in \Sigma'$ and $\mathbb{P}' A = 1$, then $\mathbb{Q} A = 1$. Thus, since $\mathcal{P}' \models X$, it follows that $\mathcal{Q} \models X$. By supposition, then, we have $\mathcal{Q} \models \psi$. Since $\psi_\Omega \in \Sigma'$, this gives $\mathbb{Q} \psi_\Omega = 1$. In other words, $\mathbb{P}'(\psi_\Omega | \varphi_\Omega) = 1$. By reversing the roles of φ and $\neg\varphi$, this same argument yields $\mathbb{P}'(\psi_\Omega | \varphi_\Omega^c) = 1$. Therefore,

$$\mathbb{P}' \psi_\Omega = \mathbb{P}' \varphi_\Omega \mathbb{P}'(\psi_\Omega | \varphi_\Omega) + \mathbb{P}' \varphi_\Omega^c \mathbb{P}'(\psi_\Omega | \varphi_\Omega^c) = \mathbb{P}' \varphi_\Omega + \mathbb{P}' \varphi_\Omega^c = 1,$$

which contradicts the fact that $\mathcal{P}' \not\models \psi$. \square

(C:soundness) Corollary 4.1.11. *If $X \subseteq \mathcal{F}$ is satisfiable, then X is consistent.*

Proof. Suppose X is inconsistent. Then $X \vdash \perp$. By Theorem 4.1.10, we have $X \models \perp$. But $\perp_\Omega = \emptyset$, so $\mathcal{P} \not\models \perp$ for all \mathcal{P} . Hence, X is not satisfiable. \square

(P:ThP-is-theory) Proposition 4.1.12. *If \mathcal{P} is a model, then $Th \mathcal{P}$ is a consistent deductive theory.*

Proof. Let $T = Th \mathcal{P}$ and suppose $T \vdash \varphi$. By Theorem 4.1.10, we have $T \models \varphi$. Since $\mathcal{P} \models T$, this implies $\mathcal{P} \models \varphi$. Hence, $\varphi \in T$, so that T is a deductive theory. Since $\emptyset = \perp_\Omega$, we have $\mathcal{P} \not\models \perp$, so that $\perp \notin T$ and T is consistent. \square

4.1.5 Karp's completeness theorem

In this subsection, we establish that our logical system is complete, meaning that every consequence of X is derivable from X . Completeness is the converse of soundness. Together, they show that the derivability and consequence relations are identical.

In Theorem 3.1.10, we showed that \vdash is σ -compact. In Theorem 4.1.17 below, we will show σ -compactness for \models , and then use this to establish completeness in Theorem 4.1.19.

It is well-known that both σ -compactness and completeness fail when we adopt the classical semantic notion of the strict model (see Example 4.4.5). In that case, only a weaker version of completeness is available. This weaker version was proven by Karp in [16]. We present Karp's version below, and then use it to establish the full completeness theorem for our notion of deductive satisfiability.

(T:Karp-compl) **Theorem 4.1.13 (Karp's completeness theorem).** *For all formulas $\varphi \in \mathcal{F}$, we have $\vdash \varphi$ if and only if $\models \varphi$.*

Proof. The only if direction is a consequence of Theorem 4.1.10. For the if direction, we appeal to Karp's completeness theorem. In [16, Theorem 5.3.2], Karp proved that $\vdash' \varphi$ if and only if $\omega \models \varphi$ for all strict models ω , where \vdash' is a certain Hilbert-type system of deduction. As noted previously, $\models \varphi$ if and only if $\omega \models \varphi$ for all strict models ω . We therefore have that $\models \varphi$ if and only if $\vdash' \varphi$. To complete the proof, we must verify that $\vdash' \varphi$ implies $\vdash \varphi$.

To accomplish this, we must first describe the differences between \vdash' and \vdash . In Karp's system, \rightarrow is a primitive symbol; for us, it is defined shorthand. This, however, causes no difficulties, since $(\varphi \rightarrow \psi) \leftrightarrow (\neg\varphi \vee \psi)$ is a tautology in Karp's system.

Recall from Theorem 3.1.17 that $\vdash \varphi$ if and only if there is a proof of φ from the axioms Λ . Karp's \vdash' differs from our \vdash only in the choice of the axioms; the notion of proof is the same. Aside from the aforementioned use of \rightarrow , this is the only difference between \vdash' and \vdash . Hence, we need only verify that each of Karp's axioms can be proven in \vdash . The axioms of Karp that are not already accounted for in Λ are:

$$(\Lambda 4) \quad \varphi \rightarrow \psi \rightarrow \varphi$$

$$(\Lambda 5) \quad (\neg\varphi \rightarrow \neg\psi) \rightarrow \psi \rightarrow \varphi$$

$$(\Lambda 6) \quad \bigwedge_{\varphi \in \Phi} (\psi \rightarrow \varphi) \rightarrow \psi \rightarrow \bigwedge \Phi$$

By Theorem 3.1.17, it suffices to prove these by natural deduction, which is entirely straightforward. \square

(R:impl-subset) **Remark 4.1.14.** As a consequence of Karp's completeness theorem, we have that φ is a tautology if and only if $\omega \models \varphi$ for all strict models ω . Hence, in any model \mathcal{P} , we have $\varphi \vdash \psi$ implies $\varphi_\Omega \subseteq \psi_\Omega$, and $\varphi \equiv \psi$ implies $\varphi_\Omega = \psi_\Omega$. If $\Omega = \mathcal{B}^{PV}$ is the set of all strict models, then both of these implications are biconditional.

With Karp's completeness theorem, we can now prove the result that was described in Remark 3.1.7.

(P:finitary-vs-infinitary)

Proposition 4.1.15. *Let $X \subseteq \mathcal{F}_{\text{fin}}$ and $\varphi \in \mathcal{F}_{\text{fin}}$. If $X \vdash \varphi$, then $X \vdash_{\text{fin}} \varphi$.*

Proof. Let $X \subseteq \mathcal{F}_{\text{fin}}$ and $\varphi \in \mathcal{F}_{\text{fin}}$. Suppose $X \vdash \varphi$. The well-known completeness theorem from finitary propositional logic states that $X \vdash_{\text{fin}} \varphi$ if and only if $\omega \models X$ implies $\omega \models \varphi$ for all strict models ω . (See, for instance, [28, Theorem 1.4.6]).

Let ω be a strict model and assume that $\omega \models X$. By Proposition 3.1.14, we may choose countable $X_0 \subseteq X$ such that $\vdash \bigwedge X_0 \rightarrow \varphi$. By Karp's completeness theorem, $\vdash \bigwedge X_0 \rightarrow \varphi$. Hence, $\omega \models \bigwedge X_0 \rightarrow \varphi$. But $\omega \models X \supseteq X_0$. Therefore, $\omega \models \varphi$. \square

4.1.6 Inductive theories and Dynkin systems

(S:ind-th-Dynk)

We briefly pause our development to make an observation about Dynkin systems. Let P be an inductive theory and fix $X \in \text{ante } P$. In Section 3.2.2, we noted that the domain of $P(\cdot | X)$ need not be closed under conjunctions and disjunctions. We are now in a position to say something in the positive direction about the structure of this set of formulas.

Let $\Omega = \mathcal{B}^{PV}$ and define

$$\Delta = \Delta(P, X) = \{\varphi_{\Omega} \mid P(\varphi | X) \text{ exists}\}. \quad (4.1.2) \quad \boxed{\text{Del-P-X}}$$

Let $A \in \mathcal{B}^{PV} = \{\varphi_{\Omega} \mid \varphi \in \mathcal{F}\}$ and choose $\varphi \in \mathcal{F}$ such that $A = \varphi_{\Omega}$. By the above definition, if $P(\varphi | X)$ exists, then $A \in \Delta$. Conversely, if $A \in \Delta$, then Remark 4.1.14 and the rule of logical equivalence imply that $P(\varphi | X)$ exists. Hence, Δ is an embedding of the domain of $P(\cdot | X)$ into \mathcal{B}^{PV} . The structure of this domain, therefore, can be understood by looking at the structure of Δ .

(P:Del-P-X)

Proposition 4.1.16. *If P be an inductive theory, then $\Delta(P, X)$ is a Dynkin system for every $X \in \text{ante } P$.*

Proof. Let P be an inductive theory and $X \in \text{ante } P$. Let $\Delta = \Delta(P, X)$ be defined as above. By the rule of logical implication, $P(\top | X) = 1$. Hence, $\Omega = \top_{\Omega} \in \Delta$, and Δ is nonempty. Let $A \in \Delta$. Choose $\varphi \in \mathcal{F}$ such that $A = \varphi_{\Omega}$. Then $A^c = (\neg\varphi)_{\Omega}$. By Corollary 3.2.7, we have $A^c \in \Delta$. Now suppose $\{A_n\} \subseteq \Delta$ is pairwise disjoint. Choose $\varphi_n \in \mathcal{F}$ such that $A_n = (\varphi_n)_{\Omega}$. For $i \neq j$, we have $\perp_{\Omega} = \emptyset = A_i \cap A_j = (\varphi_i \wedge \varphi_j)_{\Omega}$. Hence, from Remark 4.1.14, it follows that $\varphi_i \wedge \varphi_j \equiv \perp$. By the rule of logical equivalence, $P(\varphi_i \wedge \varphi_j | X) = 0$. Therefore, Theorem 3.2.24 implies $P(\bigvee_n \varphi_n | X)$ exists. But $(\bigvee_n \varphi_n)_{\Omega} = \bigcup_n A_n$, so $\bigcup_n A_n \in \Delta$, and Δ is a Dynkin system. \square

4.1.7 The full completeness theorem

(T:compactness)

Theorem 4.1.17 (σ -compactness). *A set $X \subseteq \mathcal{F}$ is satisfiable if and only if every countable subset of X is satisfiable.*

Proof. The only if part is trivial. Suppose every countable subset of X is satisfiable. Assume X is inconsistent. Then $X \vdash \perp$. By Theorem 3.1.10, there exists countable $X_0 \subseteq X$ such that $X_0 \vdash \perp$, implying that X_0 is inconsistent. By Corollary 4.1.11, we have that X_0 is not satisfiable, a contradiction. Hence, X is consistent.

Let Ω be the set of all strict models. Let

$$\Sigma = \{\varphi_\Omega \mid \varphi \in T(X) \text{ or } \neg\varphi \in T(X)\}.$$

Then Σ is a σ -algebra. If $A \in \Sigma$, choose φ such that $A = \varphi_\Omega$ and define $\mathbb{P}A = 1$ if $\varphi \in T(X)$ and 0 otherwise. By Remark 4.1.14, the function \mathbb{P} is well-defined.

Since X is consistent, $\perp \notin T(X)$. Thus, $\mathbb{P}\emptyset = \mathbb{P}\perp_\Omega = 0$. Conversely, $\top \in T(X)$, so $\mathbb{P}\Omega = \mathbb{P}\top_\Omega = 1$.

Now let $\{A_n\}_{n \in \mathbb{N}} \subseteq \Sigma$ be pairwise disjoint, and define $A = \bigcup_n A_n$. For each n , choose φ_n such that $A_n = (\varphi_n)_\Omega$, and define $\varphi = \bigvee_n \varphi_n$. Note that $A = \varphi_\Omega$. Suppose $m \neq n$. Since

$$(\varphi_m \wedge \varphi_n)_\Omega = A_m \cap A_n = \emptyset = \perp_\Omega,$$

we have $\varphi_m \wedge \varphi_n \equiv \perp$, implying that $\varphi_m \wedge \varphi_n \notin T(X)$. Since $T(X)$ is closed under conjunctions, either $\varphi_m \notin T(X)$ or $\varphi_n \notin T(X)$. This implies that there is at most one $n \in \mathbb{N}$ with $\mathbb{P}A_n = 1$. Therefore, $\sum_n \mathbb{P}A_n \in \{0, 1\}$ and

$$\begin{aligned} \sum \mathbb{P}A_n = 1 & \text{ iff there exists } n \text{ such that } \mathbb{P}A_n = 1 \\ & \text{ iff there exists } n \text{ such that } \varphi_n \in T(X) \\ & \text{ iff } \varphi \in T(X) \\ & \text{ iff } \mathbb{P}\varphi_\Omega = \mathbb{P}A = 1, \end{aligned}$$

showing that \mathbb{P} is countably additive. Thus, \mathbb{P} is a measure on (Ω, Σ) with $\mathbb{P}\Omega = 1$, and so $\mathcal{P} = (\Omega, \Sigma, \mathbb{P})$ is a model.

Now let $\varphi \in X \subseteq T(X)$ be arbitrary. Then $\varphi_\Omega \in \Sigma$, and since $\varphi \in T(X)$, we have $\mathbb{P}\varphi_\Omega = 1$, showing that $\mathcal{P} \models X$, so that X is satisfiable. \square

$\langle \text{C:compactness} \rangle$ **Corollary 4.1.18.** *A set $X \subseteq \mathcal{F}$ is satisfiable if and only if X is consistent.*

Proof. The only if part is Corollary 4.1.11. Suppose X is not satisfiable. By Theorem 4.1.17, there exists a countable subset $X_0 \subseteq X$ that is not satisfiable. By Proposition 4.1.6, the set X_0 is not strictly satisfiable. Thus, $\omega \models \neg \bigwedge X_0$ for all strict models ω , which implies $\models \neg \bigwedge X_0$. By Theorem 4.1.13, we have $\vdash \neg \bigwedge X_0$. Thus, $X \vdash \bigwedge X_0$ and $X \vdash \neg \bigwedge X_0$, showing that X is inconsistent. \square

$\langle \text{T:completeness} \rangle$ **Theorem 4.1.19 (Deductive completeness).** *For $X \subseteq \mathcal{F}$ and $\varphi \in \mathcal{F}$, we have $X \models \varphi$ if and only if $X \vdash \varphi$.*

Proof. The if part is Theorem 4.1.10. Suppose $X \not\models \varphi$. Then $X \cup \{\neg\varphi\}$ is consistent, by Theorem 3.1.13. Thus, $X \cup \{\neg\varphi\}$ is satisfiable, by Corollary 4.1.18. Let \mathcal{P} be a model with $\mathcal{P} \models X \cup \{\neg\varphi\}$. Then \mathcal{P} is an example of a model with $\mathcal{P} \models X$ and $\mathcal{P} \not\models \varphi$. Thus, $X \not\models \varphi$. \square

4.2 Inductive semantics

(S:ind-sem)

4.2.1 Inductive satisfiability

We now define a notion of satisfiability for inductive statements. Let $\mathcal{P} = (\Omega, \Sigma, \mathbb{P})$ be a model with completion $\overline{\mathcal{P}} = (\Omega, \overline{\Sigma}, \overline{\mathbb{P}})$. Let (X, φ, p) be an inductive statement. We say that \mathcal{P} *satisfies* (X, φ, p) , written $\mathcal{P} \models (X, \varphi, p)$, if there exists $Y \subseteq Th \mathcal{P}$ and $\psi \in \mathcal{F}$ such that $X \equiv Y \cup \{\psi\}$ and

$$\frac{\overline{\mathbb{P}} \varphi_{\Omega} \cap \psi_{\Omega}}{\overline{\mathbb{P}} \psi_{\Omega}} = p. \quad (4.2.1) \quad \boxed{\text{cond-prob}}$$

Note that $\mathcal{P} \models (X, \varphi, p)$ if and only if $\overline{\mathcal{P}} \models (X, \varphi, p)$. Hence, in many circumstances, we may assume without loss of generality that our models are complete.

For $Q \subseteq \mathcal{F}^{\text{IS}}$, we write $\mathcal{P} \models Q$ to mean $\mathcal{P} \models (X, \varphi, p)$ for all $(X, \varphi, p) \in Q$. A set Q is *satisfiable* if there exists a model \mathcal{P} such that $\mathcal{P} \models Q$.

The next result shows that if $\mathcal{P} \models (X, \varphi, p)$, then (4.2.1) will hold, regardless of how we decompose X into Y and ψ . As a corollary, we see that for fixed X and φ , there can be only one p such that $\mathcal{P} \models (X, \varphi, p)$.

(P:model-func)

Proposition 4.2.1. *If $X \equiv Y \cup \{\psi\} \equiv Y' \cup \{\psi'\}$ and $\mathcal{P} \models Y, Y'$, then $\psi_{\Omega} = \psi'_{\Omega}$ a.s. In particular, if $\mathcal{P} \models (X, \varphi, p)$ and $X \equiv Y \cup \{\psi\}$, where $\mathcal{P} \models Y$, then (4.2.1) holds.*

Proof. Suppose $X \equiv Y \cup \{\psi\} \equiv Y' \cup \{\psi'\}$ and $\mathcal{P} \models Y, Y'$. Using Theorem 4.1.19, we have $Y, \psi \models \psi'$, which implies $Y \models \psi \rightarrow \psi'$. Hence, $\mathcal{P} \models \psi \rightarrow \psi'$, so that $\overline{\mathbb{P}} \psi_{\Omega} \cap (\psi'_{\Omega})^c = 0$. Similarly, $\overline{\mathbb{P}} \psi'_{\Omega} \cap \psi_{\Omega}^c = 0$. Thus, $\overline{\mathbb{P}} \psi_{\Omega} \Delta \psi'_{\Omega} = 0$.

Now suppose $\mathcal{P} \models (X, \varphi, p)$ and $X \equiv Y \cup \{\psi\}$, where $\mathcal{P} \models Y$. Choose $Y' \subseteq \mathcal{F}$ and $\psi' \in \mathcal{F}$ such that $X \equiv Y' \cup \{\psi'\}$, $\mathcal{P} \models Y'$, and $\overline{\mathbb{P}} \varphi_{\Omega} \cap \psi'_{\Omega} / \overline{\mathbb{P}} \psi'_{\Omega} = p$. By the above, $\psi_{\Omega} = \psi'_{\Omega}$ a.s. Hence, $\overline{\mathbb{P}} \psi_{\Omega} = \overline{\mathbb{P}} \psi'_{\Omega}$ and $\overline{\mathbb{P}} \varphi_{\Omega} \cap \psi_{\Omega} = \overline{\mathbb{P}} \varphi_{\Omega} \cap \psi'_{\Omega}$, so that (4.2.1) holds. \square

(C:model-func)

Corollary 4.2.2. *Let \mathcal{P} be a model. If $\mathcal{P} \models (X, \varphi, p)$ and $\mathcal{P} \models (X, \varphi, p')$, then $p = p'$.*

Proof. Suppose $\mathcal{P} \models (X, \varphi, p')$. Write $X \equiv Y \cup \{\psi\}$, where $\mathcal{P} \models Y$ and $\overline{\mathbb{P}} \varphi_{\Omega} \cap \psi_{\Omega} / \overline{\mathbb{P}} \psi_{\Omega} = p'$. By Proposition 4.2.1, if $\mathcal{P} \models (X, \varphi, p)$, then (4.2.1) holds, and so $p = p'$. \square

4.2.2 Models determine theories

Given a model \mathcal{P} , we define

$$\mathbf{Th} \mathcal{P} = \{(X, \varphi, p) \in \mathcal{F}^{\text{IS}} \mid \mathcal{P} \models (X, \varphi, p)\}. \quad (4.2.2) \quad \boxed{\text{bTh-sP-def}}$$

$\langle \text{L:model-entire} \rangle$ **Lemma 4.2.3.** *Let $\mathcal{P} = (\Omega, \Sigma, \mathbb{P})$ be a model. Let $X, Y \subseteq \mathcal{F}$ and $\psi \in \mathcal{F}$. Assume $X \equiv Y \cup \{\psi\}$ and $\mathcal{P} \models Y$. Then $\mathcal{P} \models (X, \varphi, 1)$ if and only if $\psi_\Omega \in \Sigma$ and $\psi_\Omega \subseteq \varphi_\Omega$ a.s.*

Proof. Without loss of generality, assume \mathcal{P} is complete. Suppose $\mathcal{P} \models (X, \varphi, 1)$. Write $X \equiv Y' \cup \{\psi'\}$, where $\mathcal{P} \models Y'$ and $\mathbb{P} \varphi_\Omega \cap \psi'_\Omega / \mathbb{P} \psi'_\Omega = 1$. By Proposition 4.2.1, we have $\psi_\Omega = \psi'_\Omega$ a.s. Hence, $\mathbb{P} \varphi_\Omega \cap \psi_\Omega / \mathbb{P} \psi_\Omega = 1$, and this gives $\mathbb{P} \varphi_\Omega^c \cap \psi_\Omega = 0$. Conversely, suppose $\psi_\Omega \in \Sigma$ and $\psi_\Omega \subseteq \varphi_\Omega$ a.s. Then $\mathbb{P} \psi_\Omega \cap \varphi_\Omega^c = 0$, which implies $\mathbb{P} \psi_\Omega \cap \varphi_\Omega = \mathbb{P} \psi_\Omega$, so that $\mathcal{P} \models (X, \varphi, 1)$. \square

$\langle \text{T:model-ind-th} \rangle$ **Theorem 4.2.4.** *If \mathcal{P} is a model, then $\mathbf{Th} \mathcal{P}$ is a complete inductive theory with root *Taut*.*

Proof. Let \mathcal{P} be a model and let $P = \mathbf{Th} \mathcal{P}$. Without loss of generality, assume \mathcal{P} is complete. We first show that P is admissible. Suppose $(X, \varphi, p) \in P$, $X' \equiv X$, and $\varphi' \equiv_X \varphi$. Choose Y and ψ such that $X \equiv Y \cup \{\psi\}$, $\mathcal{P} \models Y$, and $\mathbb{P} \varphi_\Omega \cap \psi_\Omega / \mathbb{P} \psi_\Omega = p$. Then $Y, \psi \vdash \varphi' \leftrightarrow \varphi$, so that $Y \vdash \psi \rightarrow (\varphi \leftrightarrow \varphi')$. But $\mathcal{P} \models Y$, so $\mathbb{P} \psi_\Omega \cap (\varphi'_\Omega \Delta \varphi_\Omega) = 0$. But

$$(\varphi'_\Omega \Delta \varphi_\Omega) \cap \psi_\Omega = (\varphi'_\Omega \cap \psi_\Omega) \Delta (\varphi_\Omega \cap \psi_\Omega).$$

Thus, $\varphi'_\Omega \cap \psi_\Omega = \varphi_\Omega \cap \psi_\Omega$ a.s. Since \mathcal{P} is complete, this gives $\varphi'_\Omega \cap \psi_\Omega \in \Sigma$ and $\mathbb{P} \varphi'_\Omega \cap \psi_\Omega / \mathbb{P} \psi_\Omega = \mathbb{P} \varphi_\Omega \cap \psi_\Omega / \mathbb{P} \psi_\Omega = p$. Since $X' \equiv X \equiv Y \cup \{\psi\}$, it follows that $(X', \varphi', p) \in P$.

Now assume $(X', \varphi', p') \in P$. By Corollary 4.2.2, we have $p = p'$, and therefore P is admissible.

We next show that P is entire. Throughout the proof of entirety, we fix $X \in \text{ante } P$. Choose $(X, \varphi', p') \in P$. Write $X \equiv Y \cup \{\eta\}$, where $\mathcal{P} \models Y$ and $\mathbb{P} \varphi'_\Omega \cap \eta_\Omega / \mathbb{P} \eta_\Omega = p'$.

Suppose $X \vdash \varphi$. Then $Y \vdash \eta \rightarrow \varphi$. Hence $\mathcal{P} \models \eta \rightarrow \varphi$, which implies $\mathbb{P} \eta_\Omega \cap \varphi_\Omega^c = 0$. Since $\eta_\Omega \in \Sigma$, it follows that $\mathbb{P} \eta_\Omega \cap \varphi_\Omega = \mathbb{P} \eta_\Omega$, and therefore $P(\varphi \mid X) = 1$. Thus, P satisfies the rule of logical implication.

Suppose $P(\psi \mid X, \varphi) = 1$. Since $X \cup \{\varphi\} \equiv Y \cup \{\eta \wedge \varphi\}$, Lemma 4.2.3 gives $\eta_\Omega \cap \varphi_\Omega \subseteq \psi_\Omega$ a.s. Thus, $\eta_\Omega = (\eta_\Omega \cap \varphi_\Omega) \cup (\eta_\Omega \cap \varphi_\Omega^c) \subseteq \psi_\Omega \cup \varphi_\Omega^c$ a.s. Since $(\varphi \rightarrow \psi)_\Omega = \psi_\Omega \cup \varphi_\Omega^c$, Lemma 4.2.3 implies $P(\varphi \rightarrow \psi \mid X) = 1$, and P satisfies the rule of material implication.

Suppose $P(\varphi \mid X) = 1$ and $\varphi \vdash \psi$. By Lemma 4.2.3, we have $\eta_\Omega \subseteq \varphi_\Omega$ a.s. Remark 4.1.14 shows that $\varphi_\Omega \subseteq \psi_\Omega$. Hence, $\eta_\Omega \subseteq \psi_\Omega$ a.s., so that Lemma 4.2.3 implies $P(\psi \mid X) = 1$. Now suppose $X' \in \text{ante } P$ and $X' \vdash X$. Write $X' = Y' \cup \{\eta'\}$, where $\mathcal{P} \models Y'$ and $\eta'_\Omega \in \Sigma$. Then $Y', \eta' \vdash Y, \eta$, so that $Y' \vdash \eta' \rightarrow \eta$. Hence, $\mathcal{P} \models \eta' \rightarrow \eta$, which gives $\mathbb{P} \eta'_\Omega \cap \eta_\Omega^c = 0$. Thus, $\eta'_\Omega \subseteq \eta_\Omega$ a.s. It follows that $\eta'_\Omega \subseteq \varphi_\Omega$ a.s., so that Lemma 4.2.3 gives $P(\varphi \mid X') = 1$, and P satisfies the rule of deductive transitivity.

Suppose $X \vdash \neg(\varphi \wedge \psi)$ and two of the probabilities in (3.2.1) exist. Then $Y \vdash \eta \rightarrow \neg(\varphi \wedge \psi)$, so that $\mathbb{P} \eta_\Omega \cap \varphi_\Omega \cap \psi_\Omega = 0$. Let $\varphi' = \varphi \wedge \eta$ and $\psi' = \psi \wedge \eta$. Then $\mathbb{P} \varphi'_\Omega \cap \psi'_\Omega = 0$. Since two of the probabilities in (3.2.1) exist, two of the

sets $\varphi'_\Omega \cup \psi'_\Omega$, φ'_Ω , and ψ'_Ω are in Σ . Since $\varphi'_\Omega \cap \psi'_\Omega$ is also in Σ and Σ is a σ -algebra, it follows that all three sets are in Σ and $\mathbb{P} \varphi'_\Omega \cup \psi'_\Omega = \mathbb{P} \varphi'_\Omega + \mathbb{P} \psi'_\Omega$. From here, (3.2.1) follows immediately, and P satisfies the addition rule.

By Proposition 4.2.1, we have that $P(\varphi \mid X)$ exists and is positive if and only if $\mathbb{P} \varphi_\Omega \cap \eta_\Omega / \mathbb{P} \eta_\Omega = p$, for some $p > 0$. Similarly, $P(\varphi \wedge \psi \mid X)$ exists and is positive if and only if $\mathbb{P} \varphi_\Omega \cap \psi_\Omega \cap \eta_\Omega / \mathbb{P} \eta_\Omega = r$, for some $r > 0$. Since $X \cup \{\varphi\} \equiv Y \cup \{\varphi \wedge \eta\}$, Proposition 4.2.1 also gives that $P(\psi \mid X, \varphi)$ exists and is positive if and only if $\mathbb{P} \varphi_\Omega \cap \psi_\Omega \cap \eta_\Omega / \mathbb{P} \varphi_\Omega \cap \eta_\Omega = q$, for some $q > 0$. From this, it follows that if two of the probabilities in (3.2.2) exist and are positive, then all three exist and are positive, and $pq = r$. Hence, P satisfies the multiplication rule.

Now suppose $P(\varphi_n \mid X) = p_n$ for all n , and $X, \varphi_n \vdash \varphi_{n+1}$. By Proposition 4.2.1, we have $\mathbb{P}(\varphi_n)_\Omega \cap \eta_\Omega / \mathbb{P} \eta_\Omega = p_n$. We also have $Y \vdash \varphi_n \wedge \eta \rightarrow \varphi_{n+1}$, so that $\mathcal{P} \models \varphi_n \wedge \eta \rightarrow \varphi_{n+1}$, which gives $\mathbb{P}(\varphi_n)_\Omega \cap \eta_\Omega \cap (\varphi_{n+1})_\Omega^c = 0$. Hence, $(\varphi_n)_\Omega \cap \eta_\Omega \subseteq (\varphi_{n+1})_\Omega$ a.s. This gives $(\varphi_n)_\Omega \cap \eta_\Omega \subseteq (\varphi_{n+1})_\Omega \cap \eta_\Omega$ a.s. Since

$$(\bigvee_n \varphi_n)_\Omega \cap \eta_\Omega = (\bigcup_n (\varphi_n)_\Omega) \cap \eta_\Omega = \bigcup_n ((\varphi_n)_\Omega \cap \eta_\Omega),$$

it follows that $(\bigvee_n \varphi_n)_\Omega \cap \eta_\Omega \in \Sigma$ and, using continuity from below, we have $\mathbb{P}(\bigvee_n \varphi_n)_\Omega \cap \eta_\Omega / \mathbb{P} \eta_\Omega = \lim_{n \rightarrow \infty} \mathbb{P}(\varphi_n)_\Omega \cap \eta_\Omega / \mathbb{P} \eta_\Omega$. Therefore, P satisfies the continuity rule, and P is entire.

We next show that P is complete. Suppose $P(\varphi \mid X)$ exists. Then we may write $X \equiv Y \cup \{\eta\}$, where $\mathcal{P} \models Y$, $\mathbb{P} \varphi_\Omega \cap \eta_\Omega$ exists, and $\mathbb{P} \eta_\Omega > 0$. Assume $P(\psi \mid X)$ also exists. From Proposition 4.2.1, it follows that $\mathbb{P} \psi_\Omega \cap \eta_\Omega$ also exists. Since Σ is a σ -algebra, we have

$$(\varphi \wedge \psi)_\Omega \cap \eta_\Omega = \varphi_\Omega \cap \psi_\Omega \cap \eta_\Omega = (\varphi_\Omega \cap \eta_\Omega) \cap (\psi_\Omega \cap \eta_\Omega) \in \Sigma.$$

Hence, $\mathbb{P}(\varphi \wedge \psi)_\Omega \cap \eta_\Omega$ exists, and so therefore, $P(\varphi \wedge \psi \mid X)$ exists, showing that P satisfies Definition 3.3.1(i).

Now suppose $X \in \text{ante } P$. Then we may write $X \equiv Y \cup \{\psi\}$, where $\mathcal{P} \models Y$ and $\mathbb{P} \psi_\Omega > 0$. Assume also that $X \cup \{\varphi\} \in \text{ante } P$. Since $X \cup \{\varphi\} \equiv Y \cup \{\varphi \wedge \psi\}$, Proposition 4.2.1 implies that $\mathbb{P} \varphi_\Omega \cap \psi_\Omega$ exists. Hence, $P(\varphi \mid X)$ exists, and so P satisfies Definition 3.3.1(ii), showing that P is complete.

Since P is complete, it is therefore semi-closed. We next show that P is closed. Assume $S \subseteq \mathcal{F}$ is nonempty with $P(\theta \mid X) = 1$ for all $\theta \in S$. Then we may write $X \equiv Y \cup \{\psi\}$, where $\psi_\Omega \subseteq \theta_\Omega$ a.s. for all $\theta \in S$. Let $S' = \{\psi \rightarrow \theta \mid \theta \in S\}$ and $Y' = Y \cup S'$. By Lemma 3.1.22, we have $X \cup S \equiv Y' \cup \{\psi\}$. Also, for any $\theta \in S$, we have $\Omega = \psi_\Omega^c \cup \psi_\Omega \subseteq \psi_\Omega^c \cup \theta_\Omega = (\psi \rightarrow \theta)_\Omega$ a.s., which gives $\mathcal{P} \models \psi \rightarrow \theta$. Hence, $\mathcal{P} \models Y'$, and therefore $P(\cdot \mid X) = P(\cdot \mid X, S)$, showing that P satisfies the rule of deductive extension and is thus closed.

Finally, we show that P is connected with root \mathbf{Taut} . Let $P_0 = P \downarrow_{\mathbf{Taut}}$. Corollary 3.4.8 implies that P_0 is a complete pre-theory with root T_0 . We will show that P_0 is a basis for P . Let $X \in \text{ante } P$. Write $X \equiv Y \cup \{\psi\}$, where $\mathcal{P} \models Y$ and $\mathbb{P} \psi_\Omega > 0$. Since $\{\psi\} = \emptyset \cup \{\psi\}$, we have $\{\psi\} \in \text{ante } P_0$. Let $\theta \in Y$ be arbitrary. Since $\mathcal{P} \models Y$, we have $\mathbb{P} \theta_\Omega = 1$. Therefore, $P(\theta \mid \psi) = 1$, which gives $P_0(\theta \mid \psi) = 1$. This shows that $Y \subseteq \tau(P_0; \{\psi\})$, proving that P_0 is a basis for P . \square

$\langle \text{R:th-of-th} \rangle$ **Remark 4.2.5.** The relationship between $\mathbf{Th} \mathcal{S}$ in (4.2.2) and $Th \mathcal{S}$ in (4.1.1) is that if $P = \mathbf{Th} \mathcal{S}$, then $T_P = Th \mathcal{S}$. To see this, let $P = \mathbf{Th} \mathcal{S}$ and note that by Proposition 3.4.5, we have $\theta \in T_P$ if and only if $P(\theta \mid \mathbf{Taut}) = 1$, which holds if and only if $\mathcal{S} \models (\mathbf{Taut}, \theta, 1)$. But this holds if and only if $\mathbb{P} \theta_\Omega = 1$, and this is the definition of $\mathcal{S} \models \theta$.

4.2.3 Theories determine models

$\langle \text{T:ind-th-model} \rangle$ **Theorem 4.2.6.** *Let P be a complete inductive theory with root T_0 . Then there exists a model \mathcal{S} such that $T_P = Th \mathcal{S}$ and $P = \mathbf{Th} \mathcal{S} \downarrow_{[T_0, Th \mathcal{S}]}$. In particular, every inductive theory is satisfiable.*

Proof. Let P be a complete inductive theory with root T_0 . Let $\Omega = \mathbf{B}^{PV}$ and let $\Sigma = \{\varphi_\Omega \mid P(\varphi \mid T_0) \text{ exists}\}$. Since $\Omega = \top_\Omega$, the rule of logical implication implies $\Omega \in \Sigma$. Since $\varphi_\Omega^c = (\neg\varphi)_\Omega$, Corollary 3.2.7 implies that Σ is closed under complements. Let $\{A_n\} \subseteq \Sigma$ be pairwise disjoint, and let $A = \bigcup_n A_n$. For each n , choose φ_n such that $A_n = (\varphi_n)_\Omega$, and note that $A = \varphi_\Omega$, where $\varphi = \bigvee_n \varphi_n$. Also note that since $\{A_n\}$ are pairwise disjoint, we have $\vdash \neg(\varphi_i \wedge \varphi_j)$ for all $1 \leq i < j < \infty$. Hence, Theorem 3.2.24 implies $A \in \Sigma$, so that Σ is closed under countable, pairwise disjoint unions. It follows that Σ is a Dynkin system. Since P is complete, Definition 3.3.1(i) implies that Σ is closed under intersections. Therefore, Σ is a σ -algebra.

By Remark 4.1.14 and the rule of logical equivalence, we may define $\mathbb{P} : \Sigma \rightarrow [0, 1]$ by $\mathbb{P} \varphi_\Omega = P(\varphi \mid T_0)$. Note that $\top_\Omega = \Omega$ and $\perp_\Omega = \emptyset$, so that $\mathbb{P} \Omega = 1$ and $\mathbb{P} \emptyset = 0$. As above, Theorem 3.2.24 implies that \mathbb{P} is countably additive, so that \mathbb{P} is a probability measure on (Ω, Σ) .

Let $\mathcal{S} = (\Omega, \Sigma, \mathbb{P})$ and let $\overline{\mathcal{S}} = (\Omega, \overline{\Sigma}, \overline{\mathbb{P}})$ be its completion. Let $A \in \overline{\Sigma} \cap \mathbf{B}^{PV}$. Since $A \in \mathbf{B}^{PV}$, we may choose $\varphi \in \mathcal{F}$ such that $A = \varphi_\Omega$. Since $A \in \overline{\Sigma}$, we may write $A = \varphi_\Omega = \psi_\Omega \cup N$, where $P(\psi \mid T_0)$ exists, $N \subseteq \eta_\Omega$, and $P(\eta \mid T_0) = 0$. By (3.2.5) and the rule of logical implication, $P(\varphi \wedge \eta \mid T_0) = 0$. Hence, $\varphi_\Omega \cap \eta_\Omega \in \Sigma$. On the other hand, $\varphi_\Omega \cap \eta_\Omega^c = \psi_\Omega \cap \eta_\Omega^c \in \Sigma$. Therefore, $A = \varphi_\Omega \in \Sigma$, and this shows that $\overline{\Sigma} \cap \mathbf{B}^{PV} \subseteq \Sigma$. Since the reverse inclusion is trivial, we have $\overline{\Sigma} \cap \mathbf{B}^{PV} = \Sigma$.

We next show that for any $\varphi, \psi \in \mathcal{F}$, we have

$$P_0(\varphi \mid T_0, \psi) = p \quad \text{iff} \quad \mathcal{S} \models (T_0 + \psi, \varphi, p). \quad (4.2.3) \quad \boxed{\text{ind-th-model}}$$

Suppose $P_0(\varphi \mid T_0, \psi) = p$. Then $P(\varphi \mid T_0, \psi) = p$. By Definition 3.3.1(ii), Lemma 3.2.10, and the multiplication rule, we have that $P(\psi \mid T_0) > 0$ and

$$\frac{P(\varphi \wedge \psi \mid T_0)}{P(\psi \mid T_0)} = p. \quad (4.2.4) \quad \boxed{\text{cond-prob-th}}$$

Hence, (4.2.1) holds, so $\mathcal{S} \models (T_0 + \psi, \varphi, p)$. Conversely, suppose that $\mathcal{S} \models (T_0 + \psi, \varphi, p)$. Then (4.2.1) holds. Since $\overline{\Sigma} \cap \mathbf{B}^{PV} = \Sigma$, the same equality holds for \mathbb{P} instead of $\overline{\mathbb{P}}$. Hence, $P(\psi \mid T_0) > 0$, $P(\varphi \wedge \psi \mid T_0)$ exists, and (4.2.4) holds, which, by the multiplication rule, gives $P(\varphi \mid T_0, \psi) = p$, proving (4.2.3).

Now, Theorem 4.2.4 implies that $\mathbf{Th} \mathcal{S}$ is a complete inductive theory with root T_{aut} , and Remark 4.2.5 implies $T(\mathbf{Th} \mathcal{S}) = Th \mathcal{S}$. By the rule of logical implication, $\mathcal{S} \models T_0$. Hence, $T_0 \in [T_{\text{aut}}, Th \mathcal{S}]$. It follows from Proposition 3.5.10 that if we define $P'_0 = \mathbf{Th} \mathcal{S} \downarrow_{T_0}$ and $P' = \mathbf{P}(P'_0)$, then $T_{P'} = Th \mathcal{S}$ and $P' = \mathbf{Th} \mathcal{S} \downarrow_{[T_0, Th \mathcal{S}]}$. By (4.2.3), we have $P_0 = P'_0$. Hence, $P = P'$, and so it follows that $T_P = Th \mathcal{S}$ and $P = \mathbf{Th} \mathcal{S} \downarrow_{[T_0, Th \mathcal{S}]}$. This proves the first claim of the theorem.

For the second claim, let P be an inductive theory with root T_0 . Let $P_0 = P \downarrow_{T_0}$, so that $P = \mathbf{P}(P_0)$. Since P_0 is a pre-theory, it is semi-closed and, therefore, has a completion. By Corollary 3.4.11, we may choose a completion \overline{P}_0 that is also a pre-theory with root T_0 . Let $\overline{P} = \mathbf{P}(\overline{P}_0)$. By Proposition 3.4.20, the set \overline{P} is a complete inductive theory with root T_0 such that $P \subseteq \overline{P}$. As shown above, we may construct a model \mathcal{S} such that $P \subseteq \overline{P} \subseteq \mathbf{Th} \mathcal{S}$. Hence, $\mathcal{S} \models P$ and P is satisfiable. \square

4.2.4 Consistency and satisfiability

Note that by deductive completeness, the notions of connectivity and strong connectivity can be completely characterized in terms of semantics. Therefore, the following theorem shows that consistency can also be characterized in terms of semantics.

$\langle \text{T:ind-satis-cons} \rangle$ **Theorem 4.2.7.** *A set $Q \subseteq \mathcal{F}^{\text{IS}}$ is consistent if and only if it is connected and satisfiable.*

Proof. Suppose Q is connected and satisfiable. Choose a model \mathcal{S} such that $\mathcal{S} \models Q$. Then $Q \subseteq \mathbf{Th} \mathcal{S}$. Theorem 4.2.4 implies $\mathbf{Th} \mathcal{S}$ is an inductive theory. Hence, Q can be extended to an inductive theory, so Q is consistent. Conversely, suppose Q is consistent. Choose an inductive theory P such that $Q \subseteq P$. By Theorem 4.2.6, there exists a model \mathcal{S} such that $\mathcal{S} \models P$. Since $Q \subseteq P$, we have $\mathcal{S} \models Q$, so Q is satisfiable. \square

Recall that $Q \vdash (X, \varphi, p)$ means that Q is consistent and $\mathbf{P}_Q(\varphi \mid X) = p$. As such, the following proposition gives a semantic characterization of inductive derivability in the special case that $X \hookrightarrow T_0$, where T_0 is the root of Q . As a corollary, we find that T_Q also has a semantic characterization.

$\langle \text{P:sound-compl-pre-th} \rangle$ **Proposition 4.2.8.** *Let Q be connected and satisfiable, so that Q is also consistent. Let T_0 be the root of Q . Then the following are equivalent:*

(i) $\mathbf{P}_Q(\varphi \mid T_0, \psi) = p$,

(ii) for all models \mathcal{S} , if $\mathcal{S} \models Q$ and $\mathcal{S} \models T_0$, then $\mathcal{S} \models (T_0 + \psi, \varphi, p)$.

Proof. Suppose that $\mathbf{P}_Q(\varphi \mid T_0, \psi) = p$, and let $\mathcal{S} \models Q$ and $\mathcal{S} \models T_0$. Then $Q \subseteq \mathbf{Th} \mathcal{S}$. By Theorems 4.2.4 and 3.3.7, this gives $\mathbf{P}_Q \subseteq \mathbf{Th} \mathcal{S}$, so that $\mathcal{S} \models (T_0 + \psi, \varphi, p)$.

Conversely, suppose that for all models \mathcal{S} , if $\mathcal{S} \models Q$ and $\mathcal{S} \models T_0$, then $\mathcal{S} \models (T_0 + \psi, \varphi, p)$. Let $P_0 = \mathbf{P}_Q \downarrow_{T_0}$. Let \overline{P} be a completion of P_0 and

define $P = \mathbf{P}(\overline{P} \downarrow_{T_0})$. Since $T_0 \in \text{ante } \overline{P}$, Corollary 3.4.8 and Proposition 3.4.20 imply that P is a complete inductive theory with root T_0 . By Theorem 4.2.6, we may choose a model \mathcal{P} such that $T_P = \text{Th } \mathcal{P}$ and $P = \mathbf{Th } \mathcal{P} \downarrow_{[T_0, \text{Th } \mathcal{P}]}$. Note that $Q \subseteq \mathbf{P}_Q = \mathbf{P}(P_0) \subseteq \mathbf{P}(\overline{P} \downarrow_{T_0}) = P \subseteq \mathbf{Th } \mathcal{P}$. Hence, $\mathcal{P} \models Q$. Also $T_0 \subseteq T_P = \text{Th } \mathcal{P}$, so that $\mathcal{P} \models T_0$. Therefore, by assumption, $\mathcal{P} \models (T_0 + \psi, \varphi, p)$. But $P = \mathbf{Th } \mathcal{P} \downarrow_{[T_0, \text{Th } \mathcal{P}]}$, so this gives $P(\varphi \mid T_0, \psi) = p$. Also, $P = \mathbf{P}(\overline{P} \downarrow_{T_0})$, which implies $P \downarrow_{T_0} = \overline{P} \downarrow_{T_0}$. Hence, $\overline{P}(\varphi \mid T_0, \psi) = p$. Since \overline{P} was arbitrary, the rule of inductive extension yields $P_0(\varphi \mid T_0, \psi) = p$, and therefore $\mathbf{P}_Q(\varphi \mid T_0, \psi) = p$. \square

(C:sound-compl-pre-th) **Corollary 4.2.9.** *Let Q be connected and satisfiable, so that Q is also consistent. Let T_0 be the root of Q . Then the following are equivalent:*

(i) $\theta \in T_Q$,

(ii) for all models \mathcal{P} , if $\mathcal{P} \models Q$ and $\mathcal{P} \models T_0$, then $\mathcal{P} \models \theta$.

Proof. Assume $\theta \in T_Q$, so that $\mathbf{P}_Q(\theta \mid T_0) = 1$. Suppose $\mathcal{P} \models Q$ and $\mathcal{P} \models T_0$. Proposition 4.2.8 implies $\mathcal{P} \models (T_0, \theta, 1)$. Since $T_0 \equiv T_0 \cup \{\top\}$ and $\mathcal{P} \models T_0$, we have $\overline{\mathbb{P}}\theta_\Omega \cap \top_\Omega / \overline{\mathbb{P}}\top_\Omega = 1$. But $\top_\Omega = \Omega$, so $\overline{\mathbb{P}}\theta_\Omega = 1$, which means $\mathcal{P} \models \theta$.

Now assume that for all models \mathcal{P} , if $\mathcal{P} \models Q$ and $\mathcal{P} \models T_0$, then $\mathcal{P} \models \theta$. As above, if $\mathcal{P} \models T_0$ and $\mathcal{P} \models \theta$, then $\mathcal{P} \models (T_0, \theta, 1)$. Hence, for all models \mathcal{P} , if $\mathcal{P} \models Q$ and $\mathcal{P} \models T_0$, then $\mathcal{P} \models (T_0, \theta, 1)$. Proposition 4.2.8 implies $\mathbf{P}_Q(\theta \mid T_0) = 1$, so that by Proposition 3.4.5, we have $\theta \in T_Q$. \square

4.2.5 Inductive consequence and completeness

Having characterized T_Q in terms of semantics, we are now ready to define the inductive consequence relation.

(D:consequence) **Definition 4.2.10.** Let $Q \subseteq \mathcal{F}^{\text{IS}}$ and $(X, \varphi, p) \in \mathcal{F}^{\text{IS}}$. We say that (X, φ, p) is a *consequence* of Q , or that Q *entails* (X, φ, p) , which we denote by $Q \models (X, \varphi, p)$, if

(i) Q is connected and satisfiable,

(ii) $X \hookrightarrow [T_0, T_Q]$, where T_0 is the root of Q , and

(iii) $\mathcal{P} \models Q$ implies $\mathcal{P} \models (X, \varphi, p)$, for all models \mathcal{P} .

Corollary 4.2.9 shows that T_Q has an entirely semantic characterization. Hence, Definition 4.2.10 is an entirely semantic definition.

Note that $\mathcal{P} \models (\mathbf{Taut}, \top, 1)$ for all models \mathcal{P} . However, $Q \vdash (\mathbf{Taut}, \top, 1)$ only if the root of Q is \mathbf{Taut} . Hence, (ii) cannot be removed if we are to have completeness.

To simplify the verification that $Q \models (X, \varphi, p)$, we can replace (iii) with

(iii)' $\mathcal{P} \models Q$ implies $\mathcal{P} \models (X, \varphi, p)$ whenever \mathcal{P} is complete and $\mathcal{P} \models T_0$.

We show this below in Theorem 4.2.13, after proving two lemmas. Let $\mathcal{Q} = (\Omega, \Sigma, \mathbb{Q})$ be a complete model and assume that $\mathbb{Q}\zeta_\Omega > 0$ for some $\zeta \in \mathcal{F}$. Define the probability measure \mathbb{P} on (Ω, Σ) by $\mathbb{P}A = \mathbb{Q}A \cap \zeta_\Omega / \mathbb{Q}\zeta_\Omega$ and let $\mathcal{P} = (\Omega, \Sigma, \mathbb{P})$. Note that if $\mathbb{Q}A = 1$, then $\mathbb{P}A = 1$. Hence, $\mathcal{Q} \models Y$ implies $\mathcal{P} \models Y$ for all $Y \subseteq \mathcal{F}$.

(L:cond-compl) **Lemma 4.2.11.** *Let \mathcal{Q} and \mathcal{P} be as above and let $\overline{\mathcal{P}} = (\Omega, \overline{\Sigma}, \overline{\mathbb{P}})$ be the completion of \mathcal{P} . Then $A \in \overline{\Sigma}$ if and only if $A \cap \zeta_\Omega \in \Sigma$, and in this case, $\overline{\mathbb{P}}A = \mathbb{P}A \cap \zeta_\Omega$.*

Proof. Suppose $A \cap \zeta_\Omega \in \Sigma$. Since $\mathbb{P}\zeta_\Omega^c = 0$ and $A \cap \zeta_\Omega^c \subseteq \zeta_\Omega^c$, we have $A \cap \zeta_\Omega^c \in \overline{\Sigma}$. Hence, $A = (A \cap \zeta_\Omega) \cup (A \cap \zeta_\Omega^c) \in \overline{\Sigma}$ and $\overline{\mathbb{P}}A = \mathbb{P}A \cap \zeta_\Omega$. Conversely, suppose $A \in \overline{\Sigma}$. Write $A = B \cup F$, where $B \in \Sigma$, $F \subseteq N$, and $N \in \Sigma$ with $\mathbb{P}N = 0$. By the definition of \mathbb{P} , we have $\mathbb{Q}N \cap \zeta_\Omega = 0$. Now write $A \cap \zeta_\Omega = (B \cap \zeta_\Omega) \cup (F \cap \zeta_\Omega)$. Then $B \cap \zeta_\Omega \in \Sigma$. Also, since $F \cap \zeta_\Omega \subseteq N \cap \zeta_\Omega$ and \mathcal{Q} is complete, it follows that $F \cap \zeta_\Omega \in \Sigma$. Therefore, $A \cap \zeta_\Omega \in \Sigma$. Moreover, this shows that $\overline{\mathbb{P}}A = \mathbb{P}B$ and $\mathbb{P}A \cap \zeta_\Omega = \mathbb{P}B \cap \zeta_\Omega$. Since $\mathbb{P}\zeta_\Omega = 1$, we have $\mathbb{P}B = \mathbb{P}B \cap \zeta_\Omega$. Therefore, $\overline{\mathbb{P}}A = \mathbb{P}A \cap \zeta_\Omega$. \square

(L:conseq-defn-simple) **Lemma 4.2.12.** *Let \mathcal{Q} and \mathcal{P} be as above. If $\mathcal{Q} \models (X, \varphi, p)$ and $\zeta \in T(X)$, then $\mathcal{P} \models (X, \varphi, p)$.*

Proof. Suppose $\mathcal{Q} \models (X, \varphi, p)$ and $\zeta \in T(X)$. Write $X \equiv Y \cup \{\psi\}$, where $\mathcal{Q} \models Y$. Then $\mathcal{P} \models Y$ also. Since $\zeta \in T(X)$, we have $X \equiv Y \cup \{\psi, \zeta\} \equiv Y \cup \{\psi \wedge \zeta\}$. By Proposition 4.2.1 and the fact that \mathcal{Q} is complete, it follows that

$$p = \frac{\mathbb{Q}\varphi_\Omega \cap \psi_\Omega \cap \zeta_\Omega}{\mathbb{Q}\psi_\Omega \cap \zeta_\Omega} = \frac{\mathbb{P}\varphi_\Omega \cap \psi_\Omega}{\mathbb{P}\psi_\Omega}.$$

Therefore, $\mathcal{P} \models (X, \varphi, p)$. \square

(T:conseq-defn-simple) **Theorem 4.2.13.** *In Definition 4.2.10, we may replace (iii) with (iii)'.*

Proof. Clearly, Definition 4.2.10 implies (iii)'. For the converse, let $Q \subseteq \mathcal{F}^{\text{IS}}$ and $(X, \varphi, p) \in \mathcal{F}^{\text{IS}}$. Assume (i), (ii), and (iii)'. To show that (iii) holds, let \mathcal{Q} be a model and assume that $\mathcal{Q} \models Q$. We want to show that $\mathcal{Q} \models (X, \varphi, p)$. As noted below (4.2.1), we may assume without loss of generality that \mathcal{Q} is complete.

By Proposition 3.4.2, we may choose X_0 such that $T_0 = T(X_0)$ and $X_0 \in \text{ante}Q$. We may then choose $(X_0, \xi, q) \in Q$, so that $\mathcal{Q} \models (X_0, \xi, q)$. Write $X_0 \equiv Y \cup \{\zeta\}$, where $\mathcal{Q} \models Y$ and $\mathbb{Q}\zeta_\Omega > 0$. Define $\mathcal{P} = (\Omega, \Sigma, \mathbb{P})$ by $\mathbb{P}A = \mathbb{Q}A \cap \zeta_\Omega / \mathbb{Q}\zeta_\Omega$. Then $\mathcal{P} \models Y$ and $\mathcal{P} \models \zeta$. Hence, $\mathcal{P} \models X_0$, which gives $\mathcal{P} \models T_0$. Let $(X', \varphi', p') \in Q$. Then $\zeta \in T_0 \subseteq T(X')$. By Lemma 4.2.12, it follows that $\mathcal{P} \models (X', \varphi', p')$. This shows that $\mathcal{P} \models Q$. We therefore have $\overline{\mathcal{P}} \models T_0$ and $\overline{\mathcal{P}} \models Q$. By (iii)', this gives $\overline{\mathcal{P}} \models (X, \varphi, p)$, which implies $\mathcal{P} \models (X, \varphi, p)$.

By (ii), we may write $X \equiv T + \psi$, where $T \in [T_0, T_Q]$. Since $T \subseteq T_Q$ and $\mathcal{Q} \models Q$, we have $\mathcal{Q} \models T$. It remains only to show that $\mathbb{Q}\varphi_\Omega \cap \psi_\Omega / \mathbb{Q}\psi_\Omega = p$.

For this, note that $\zeta \in T_0 \subseteq T(X)$, so that $X \equiv T + \psi \wedge \zeta$. Since $\mathcal{Q} \models T$, we also have $\mathcal{P} \models T$. Hence, by Proposition 4.2.1 and Lemma 4.2.11, it follows that

$$p = \frac{\mathbb{P} \varphi_\Omega \cap \psi_\Omega \cap \zeta_\Omega}{\mathbb{P} \psi_\Omega \cap \zeta_\Omega} = \frac{\mathbb{Q} \varphi_\Omega \cap \psi_\Omega}{\mathbb{Q} \psi_\Omega}.$$

Therefore, $\mathcal{Q} \models (X, \varphi, p)$, which shows that (iii) holds. \square

Having defined the inductive consequence relation, we now show that it is identical to the inductive derivability relation.

$\langle \text{T:ind-sound-comp1} \rangle$ **Theorem 4.2.14 (Inductive soundness and completeness).** *Let $Q \subseteq \mathcal{F}^{\text{IS}}$ and $(X, \varphi, p) \in \mathcal{F}^{\text{IS}}$. Then $Q \vdash (X, \varphi, p)$ if and only if $Q \models (X, \varphi, p)$.*

Proof. Suppose $Q \vdash (X, \varphi, p)$. Then Q is consistent and $\mathbf{P}_Q(\varphi \mid X) = p$. By Remark 3.4.24, we have $X \hookrightarrow [T_0, T_Q]$, where T_0 is the root of Q . Theorem 4.2.7 implies Q is connected and satisfiable. Suppose $\mathcal{P} \models Q$. Theorems 4.2.4 and 3.3.7 implies $\mathcal{P} \models \mathbf{P}_Q$. Hence, $\mathcal{P} \models (X, \varphi, p)$.

For the converse, suppose $Q \models (X, \varphi, p)$. We need to show that $\mathbf{P}_Q(\varphi \mid X) = p$. By Definition 4.2.10(ii), we may write $T(X) = T + \psi$, where $T \in [T_0, T_Q]$. Hence, it suffices to show $\mathbf{P}_Q(\varphi \mid T_0, \psi) = p$. We will do this using Proposition 4.2.8.

Suppose $\mathcal{P} \models Q$ and $\mathcal{P} \models T_0$. Then $T_0 \subseteq \text{Th } \mathcal{P}$, so by Remark 4.2.5 and Proposition 3.5.10, if we define $P = \mathbf{Th } \mathcal{P} \downarrow_{[T_0, \text{Th } \mathcal{P}]}$, then P is an inductive theory with root T_0 , and $T_P = \text{Th } \mathcal{P}$. Corollary 4.2.9 gives $T_Q \subseteq \text{Th } \mathcal{P}$, so that $T \in [T_0, \text{Th } \mathcal{P}]$. Hence, $X \hookrightarrow [T_0, \text{Th } \mathcal{P}]$. Moreover, Definition 4.2.10(iii) implies $(X, \varphi, p) \in \mathbf{Th } \mathcal{P}$. Hence, $P(\varphi \mid X) = p$. But $T(X) = T + \psi$ and $T \in [T_0, T_P]$. Therefore, $P(\varphi \mid T_0, \psi) = p$. This implies $\mathcal{P} \models (T_0 + \psi, \varphi, p)$, so by Proposition 4.2.8, we have $\mathbf{P}_Q(\varphi \mid T_0, \psi) = p$. \square

$\langle \text{R:classic-ind-th-char} \rangle$ **Remark 4.2.15.** According to Remark 3.3.8 and Theorems 4.2.7 and 4.2.14, a connected and satisfiable set $P \subseteq \mathcal{F}^{\text{IS}}$ is an inductive theory if and only if $P \models (X, \varphi, p)$ implies $(X, \varphi, p) \in P$ for all $(X, \varphi, p) \in \mathcal{F}^{\text{IS}}$.

4.2.6 Differing roots

$\langle \text{S:diff-root} \rangle$ Proposition 3.5.10 shows that if P is an inductive theory with root T_0 , and $T'_0 \in [T_0, T_P]$, then $P' = P \downarrow_{[T'_0, T_P]}$ is an inductive theory with root T'_0 . The difference between these two theories is described entirely by the sentences in $T'_0 \setminus T_0$. In P' , these sentences are part of the root, which means they are part of every antecedent. In P , they are part of T_P , which means they have probability one under every antecedent. In either case, we are assuming such a sentence is “true,” in one sense or another. It may be tempting, then, to think that these two theories are effectively the same. In fact, for any model \mathcal{P} , if $\mathcal{P} \models P$, then $\mathcal{P} \models P'$. The converse, however, is not true. The theory with the lower root, P , has fewer models, as we illustrate below in Proposition 4.2.16. In other words, by placing a hypothesis in T_P rather than T_0 , we are making a semantically stronger assumption. Intuitively, the sentences in T_0 are only

hypothetical postulates. The inductive statements in P are all assertions about the case in which we are given T_0 . A sentence $\zeta \in T_P \setminus T_0$ has a different status. In that case, we have $P(\zeta \mid T_0) = 1$. Hence, the inductive theory P is asserting that ζ is entailed (probabilistically) by T_0 .

To illustrate this fact, let $\mathcal{Q} = (\Omega, \Sigma, \mathbb{Q})$ be a complete model and assume that $\mathbb{Q}\zeta_\Omega \in (0, 1)$ for some $\zeta \in \mathcal{F}$. Define the probability measure \mathbb{P} on (Ω, Σ) by $\mathbb{P}A = \mathbb{Q}A \cap \zeta_\Omega / \mathbb{Q}\zeta_\Omega$ and let $\mathcal{P} = (\Omega, \Sigma, \mathbb{P})$. Let $(\Omega, \Sigma, \mathbb{P})$ be the completion of $(\Omega, \Sigma, \mathbb{Q})$. Note that $\mathbf{Th} \mathcal{P}$ is an inductive theory with root \mathbf{Taut} , and $T_0 = T(\zeta) \in [\mathbf{Taut}, \mathbf{Th} \mathcal{P}]$. Let $P = \mathbf{Th} \mathcal{P} \upharpoonright_{[T_0, \mathbf{Th} \mathcal{P}]}$.

(P:drop-root-fewer-models) **Proposition 4.2.16.** *With notation as above, we have $\mathcal{Q} \models P$, but $\mathcal{Q} \not\models \mathbf{Th} \mathcal{P}$.*

Proof. Since $\mathcal{P} \models (\mathbf{Taut}, \zeta, 1)$, but $\mathbb{Q}\zeta_\Omega < 1$, we have $\mathcal{Q} \not\models \mathbf{Th} \mathcal{P}$. Suppose $P(\varphi \mid T_0, \psi) = p$. Then $\mathbb{P}\varphi_\Omega \cap \psi_\Omega / \mathbb{P}\psi_\Omega = p$. By Lemma 4.2.11,

$$p = \frac{\mathbb{P}\varphi_\Omega \cap \psi_\Omega \cap \zeta_\Omega}{\mathbb{P}\psi_\Omega \cap \zeta_\Omega} = \frac{\mathbb{Q}\varphi_\Omega \cap \psi_\Omega \cap \zeta_\Omega}{\mathbb{Q}\psi_\Omega \cap \zeta_\Omega} = \frac{\mathbb{Q}\varphi_\Omega \cap (\psi \wedge \zeta)_\Omega}{\mathbb{Q}(\psi \wedge \zeta)_\Omega}.$$

Since $T_0 \cup \{\psi\} \equiv \emptyset \cup \{\zeta \wedge \psi\}$ and $\mathcal{Q} \models \emptyset$, we have $\mathcal{Q} \models (T_0 \cup \{\psi\}, \varphi, p)$. This shows that $\mathcal{Q} \models P \upharpoonright_{T_0}$. By Proposition 3.5.10, it follows that $\mathcal{Q} \models P$. \square

4.2.7 The semantics of inductive conditions

Remark 4.2.15 gives an entirely semantic characterization of inductive theories. Consequently, inductive conditions can also be characterized semantically. We are therefore ready to extend the notions of satisfiability and consequence to inductive conditions.

We say that a model \mathcal{P} *satisfies* an inductive condition \mathcal{C} if $\mathcal{P} \models P$ for some $P \in \mathcal{C}$. An inductive condition is *satisfiable* if $\mathcal{P} \models \mathcal{C}$ for some model \mathcal{P} . The following proposition shows that this notion of satisfiability is an extension of our previous definition.

(P:set-cond-satis) **Proposition 4.2.17.** *Let $Q \subseteq \mathcal{F}^{\text{IS}}$ be connected and let \mathcal{P} be a model. Then $\mathcal{P} \models Q$ if and only if $\mathcal{P} \models \mathcal{C}_Q$.*

Proof. Suppose $\mathcal{P} \models \mathcal{C}_Q$. Choose $P \in \mathcal{C}_Q$ such that $\mathcal{P} \models P$. Since $Q \subseteq P$, we have $\mathcal{P} \models Q$. Conversely, suppose $\mathcal{P} \models Q$. Theorem 4.2.7 implies Q is consistent, so we may define \mathbf{P}_Q . Since $Q \subseteq \mathbf{Th} \mathcal{P}$ and Theorem 4.2.4 implies $\mathbf{Th} \mathcal{P}$ is an inductive theory, we have $\mathbf{P}_Q \subseteq \mathbf{Th} \mathcal{P}$. That is, $\mathcal{P} \models \mathbf{P}_Q$. But $\mathbf{P}_Q \in \mathcal{C}_Q$. Hence, $\mathcal{P} \models \mathcal{C}_Q$. \square

The next two results show that both the consistency of \mathcal{C} and the deductive theory $T_{\mathcal{C}}$ have semantic characterizations.

(T:ind-satis-cons-IC) **Theorem 4.2.18.** *An inductive condition is consistent if and only if it is satisfiable.*

Proof. Let \mathcal{C} be an inductive condition. If \mathcal{C} is satisfiable, then by definition, it is nonempty, and therefore consistent. For the converse, suppose \mathcal{C} is consistent. Choose $P \in \mathcal{C}$. By Theorem 4.2.6, we may choose a model \mathcal{P} such that $\mathcal{P} \models P$. Hence, $\mathcal{P} \models \mathcal{C}$, and \mathcal{C} is satisfiable. \square

Proposition 4.2.19. *Let \mathcal{C} be a satisfiable inductive condition, so that \mathcal{C} is also consistent. Let T_0 be the root of \mathcal{C} . Then the following are equivalent:*

- (i) $\theta \in T_{\mathcal{C}}$,
- (ii) for all models \mathcal{P} , if $\mathcal{P} \models \mathcal{C}$ and $\mathcal{P} \models T_0$, then $\mathcal{P} \models \theta$.

Proof. Suppose $\theta \in T_{\mathcal{C}}$. Assume $\mathcal{P} \models \mathcal{C}$ and $\mathcal{P} \models T_0$. Choose $P \in \mathcal{C}$ such that $\mathcal{P} \models P$. Since $T_{\mathcal{C}} = \bigcap \{T_P \mid P \in \mathcal{C}\}$, we have $\theta \in T_P$. From Corollary 4.2.9, it follows that $\mathcal{P} \models \theta$.

Now suppose that for all models \mathcal{P} , if $\mathcal{P} \models \mathcal{C}$ and $\mathcal{P} \models T_0$, then $\mathcal{P} \models \theta$. Assume $\theta \notin T_{\mathcal{C}}$. Then we may choose $P \in \mathcal{C}$ such that $\theta \notin T_P$. By Corollary 4.2.9, we may choose a model \mathcal{P} such that $\mathcal{P} \models P$, $\mathcal{P} \models T_0$, and $\mathcal{P} \not\models \theta$. Since $\mathcal{P} \models P$ and $P \in \mathcal{C}$, we have $\mathcal{P} \models \mathcal{C}$. Hence, by our initial supposition, $\mathcal{P} \models \theta$, a contradiction. \square

Having characterized $T_{\mathcal{C}}$ in terms of semantics, we can now extend the consequence relation to inductive conditions.

(D:consequence-IC) **Definition 4.2.20.** We say that $(X, \varphi, p) \in \mathcal{F}^{\text{IS}}$ is a *consequence* of an inductive condition \mathcal{C} , or that \mathcal{C} *entails* (X, φ, p) , which we denote by $\mathcal{C} \models (X, \varphi, p)$, if

- (i) \mathcal{C} is satisfiable,
- (ii) $X \leftrightarrow [T_0, T_{\mathcal{C}}]$, where T_0 is the root of \mathcal{C} , and
- (iii) $\mathcal{P} \models \mathcal{C}$ implies $\mathcal{P} \models (X, \varphi, p)$, for all models \mathcal{P} .

From Proposition 4.2.17, it follows that for any connected $Q \subseteq \mathcal{F}^{\text{IS}}$, we have $Q \models (X, \varphi, p)$ if and only if $\mathcal{C}_Q \models (X, \varphi, p)$. Hence, Definition 4.2.20 is a generalization of Definition 4.2.10.

As with Definition 4.2.10, we cannot remove (ii). In fact, in Section 4.4.1, we provide an example where (i) and (iii) above are satisfied, but (ii) fails because X is too large. (See Remark 4.4.4.)

(T:ind-sound-compl-IC) **Theorem 4.2.21 (Soundness and completeness for conditions).** *Let \mathcal{C} be an inductive condition and $(X, \varphi, p) \in \mathcal{F}^{\text{IS}}$. Then $\mathcal{C} \vdash (X, \varphi, p)$ if and only if $\mathcal{C} \models (X, \varphi, p)$.*

Proof. Suppose $\mathcal{C} \vdash (X, \varphi, p)$. Then \mathcal{C} is consistent and $\mathbf{P}_{\mathcal{C}}(\varphi \mid X) = p$. By Remark 3.4.24, we have $X \leftrightarrow [T_0, T_{\mathcal{C}}]$, where T_0 is the root of \mathcal{C} . Theorem 4.2.18 implies \mathcal{C} is satisfiable. Suppose $\mathcal{P} \models \mathcal{C}$. Choose $P \in \mathcal{C}$ such that $\mathcal{P} \models P$. Since $\mathbf{P}_{\mathcal{C}} \subseteq P$, we have $P(\varphi \mid X) = p$. By Remark 4.2.15, this implies $P \models (X, \varphi, p)$. But $T_{\mathcal{C}} \subseteq T_P$, so $X \leftrightarrow [T_0, T_P]$. Therefore, Definition 4.2.10(iii) gives $\mathcal{P} \models (X, \varphi, p)$.

For the converse, suppose $\mathcal{C} \models (X, \varphi, p)$. We need to show that $\mathbf{P}_{\mathcal{C}}(\varphi \mid X) = p$. By Definition 4.2.20(ii), we may write $T(X) = T + \psi$, where $T \in [T_0, T_{\mathcal{C}}]$. Hence, it suffices to show that $\mathbf{P}_{\mathcal{C}}(\varphi \mid T_0, \psi) = p$.

Let $P \in \mathcal{C}$ be arbitrary. If $\mathcal{P} \models P$, then $\mathcal{P} \models \mathcal{C}$, so by supposition, $\mathcal{P} \models (X, \varphi, p)$. Since $T_{\mathcal{C}} \subseteq T_P$, we have $X \leftrightarrow [T_0, T_P]$. Thus, $P \models (X, \varphi, p)$.

By Remark 4.2.15, this gives $P(\varphi \mid X) = p$. But $T(X) = T + \psi$, where $T \in [T_0, T_C] \subseteq [T_0, T_P]$. Therefore, $P(\varphi \mid T_0, \psi) = p$. Since P was arbitrary, it follows that $(T_0 + \psi, \varphi, p) \in \bigcap \mathcal{C}^0 \subseteq \mathbf{P}_C$, so $\mathbf{P}_C(\varphi \mid T_0, \psi) = p$. \square

4.3 Counterexamples and resolutions I

(S:Examples1) In this section, we present several examples that serve to illustrate the necessity of rules (R8) and (R9). More specifically, entire sets, which are closed under only (R1)–(R7), exhibit a number of pathological behaviors. These behaviors were alluded to in Chapter 3. In this section, we provide concrete examples.

To develop these examples, we must expand our tools for creating inductive theories and entire sets. We do this in Sections 4.3.1 and 4.3.2.

4.3.1 Every probability space is a model

(S:every-prob-sp) A model is a particular type of probability space, namely one in which Ω is a set of strict models. Theorem 4.3.1 below shows that every probability space, regardless of Ω , is isomorphic to a model. More specifically, for every probability space, there is an appropriate choice of PV such that the given probability space is isomorphic to a model in $\mathcal{F}(PV)$.

Later, in Theorem 5.4.2, we will give a version of this result in predicate logic. Theorem 5.4.2 will not only be concerned with an arbitrary probability space, but also with an arbitrary collection of random variables on that space.

(T:prob-sp-model-iso) **Theorem 4.3.1.** *Let PV be a given set of propositional variables and let $\mathcal{F} = \mathcal{F}(PV)$. Let (S, Γ, ν) be an arbitrary probability space. Then (S, Γ, ν) has a subspace that is isomorphic to a model in \mathcal{F} . If $\text{card}(PV) \geq \text{card}(\Gamma)$, then (S, Γ, ν) itself is isomorphic to a model. In particular, every probability space is isomorphic to a model in $\mathcal{F}(PV)$ for an appropriate choice of PV .*

Proof. Let $G : PV \rightarrow \Gamma$. If $\text{card}(PV) \geq \text{card}(\Gamma)$, then take G to be surjective. Extend G to \mathcal{F} by $G\neg\varphi = (G\varphi)^c$ and $G\bigwedge_n \varphi_n = \bigcap_n G\varphi_n$. Let $\Theta = G\mathcal{F} \subseteq \Gamma$ and note that Θ is a σ -algebra. Hence, $(S, \Theta, \nu|_{\Theta})$ is a measure subspace of (S, Γ, ν) . If $\text{card}(PV) \geq \text{card}(\Gamma)$, then $\Theta = \Gamma$. We abuse notation and simply write ν for both ν and $\nu|_{\Theta}$. We will show that (S, Θ, ν) is isomorphic to a model in \mathcal{F} . More specifically, if $\Omega = \mathbf{B}^{PV}$ is the set of all strict models, and $\Sigma = \mathcal{B}^{BV} = \{\varphi_{\Omega} \mid \varphi \in \mathcal{F}\}$, then we will construct a probability measure \mathbb{P} on (Ω, Σ) such that (S, Γ, ν) and $(\Omega, \Sigma, \mathbb{P})$ are isomorphic.

For $x \in S$, define the strict model ω^x by $\omega^x \models \mathbf{r}$ if and only if $x \in G\mathbf{r}$, for all $\mathbf{r} \in PV$. By formula induction, it follows that $\omega^x \models \varphi$ if and only if $x \in G\varphi$, for all $\varphi \in \mathcal{F}$. Define $h : S \rightarrow \Omega$ by $hx = \omega^x$. After constructing \mathbb{P} , we will show that h induces an isomorphism.

To construct \mathbb{P} , we first prove that $\varphi \equiv \psi$ implies $G\varphi = G\psi$. Recall the set of axioms, Λ , defined in Section 3.1.3. It is straightforward to verify that $G\varphi = S$ if $\varphi \in \Lambda$ is an axiom. Suppose φ is a tautology. Then there is a proof of φ from \emptyset . Using induction of the length of the proof, as in the proof of Proposition

3.1.16, one readily verifies that $G\varphi = S$. Now suppose $\varphi \equiv \psi$. Then $\varphi \leftrightarrow \psi$ is a tautology. Hence $G\varphi \leftrightarrow \psi = S$. But $G\varphi \leftrightarrow \psi = (G\varphi \Delta G\psi)^c$. Therefore, $G\varphi = G\psi$, proving the claim.

By Remark 4.1.14, we have $\varphi_\Omega = \psi_\Omega$ if and only if $\varphi \equiv \psi$. Hence, if $\varphi_\Omega = \psi_\Omega$, then $G\varphi = G\psi$. We may therefore define $g : \Sigma \rightarrow \Theta$ by $g\varphi_\Omega = G\varphi$. Let $A \in \Sigma$ and choose $\varphi \in \mathcal{F}$ such that $A = \varphi_\Omega = \{\omega \in \Omega \mid \omega \models \varphi\}$. Then

$$\begin{aligned} x \in h^{-1}A & \text{ iff } hx \in A \\ & \text{ iff } \omega^x \in \varphi_\Omega \\ & \text{ iff } \omega^x \models \varphi \\ & \text{ iff } x \in G\varphi = g\varphi_\Omega = gA. \end{aligned}$$

Hence, $h^{-1} = g$, so that $h^{-1}\varphi_\Omega = G\varphi$. In particular, this shows that h is (Θ, Σ) -measurable, so we may define $\mathbb{P} = \nu \circ h^{-1} = \nu \circ g$, making $\mathcal{P} = (\Omega, \Sigma, \mathbb{P})$ a model.

To verify that h induces an isomorphism from (S, Θ, ν) to \mathcal{P} , we must check that for each $U \in \Theta$, there exists $A \in \Sigma$ such that $h^{-1}A = U$ ν -a.s. Let $U \in \Theta$ and choose $\varphi \in \mathcal{F}$ such that $U = G\varphi$. Then $A = \varphi_\Omega \in \Sigma$ and, by the above, $h^{-1}A = G\varphi = U$. \square

4.3.2 Dynkin spaces

(S:Dynkin-sp) Theorem 4.2.4 gives us the means to construct inductive theories. For the first set of examples in the section, however, we need to construct entire sets that are not inductive theories. We will do this using Dynkin systems. More specifically, we will define what we call Dynkin spaces, a generalization of probability spaces that use Dynkin systems instead of σ -algebras. Then, analogous to Theorem 4.2.4, we will use these Dynkin spaces to construct entire sets.

(D:Dynk-sp) **Definition 4.3.2.** A *Dynkin space* is a triple, (S, Δ, ρ) , where S is a nonempty set, Δ is a Dynkin system on S , and $\rho : \Delta \rightarrow [0, 1]$ satisfies

- (i) $\rho\Omega = 1$,
- (ii) if $A, B \in \Delta$ with $A \subseteq B$, then $\rho B \setminus A = \rho B - \rho A$, and
- (iii) if $\{A_n\} \subseteq \Delta$ with $A_n \subseteq A_{n+1}$, then $\rho \bigcup A_n = \lim \rho A_n$.

Let (S, Δ, ρ) be a Dynkin space. A set $A \in \Delta$ is called *(Dynkin) measurable*. A set $A \in \Delta$ is *null* if $A \in \Delta$ and $\rho A = 0$. Note that a measurable subset of a null set is a null set. A Dynkin space is *complete* if every subset of a null set is a null set.

Proposition 4.3.3. *If (S, Δ, ρ) is a Dynkin space, then ρ is countably additive. That is, if $\{A_n\} \subseteq \Delta$ is pairwise disjoint, then $\rho \bigcup A_n = \sum \rho A_n$.*

Proof. Let $A, B \in \Delta$ be disjoint. Since Δ is a Dynkin system, it is closed under countable, pairwise disjoint unions. Hence, $A \cup B \in \Delta$. By Definition 4.3.2(ii), we have $\rho B = \rho A \cup B - \rho A$. Therefore, ρ is finitely additive. From Definition 4.3.2(iii), we obtain $\rho \bigcup A_n = \lim_n \sum_1^n \rho A_j$. \square

(P:Dynk-complete-reverse) **Proposition 4.3.4.** *Let (S, Δ, ρ) be a complete Dynkin space. If $A \in \Delta$, $\rho A = 1$, and $A \subseteq B$, then $\rho B = 1$.*

Proof. Suppose $A \in \Delta$, $\rho A = 1$, and $A \subseteq B$. Then $A^c \in \Delta$, $\rho A^c = 0$, and $B^c \subseteq A^c$. Since (S, Δ, ρ) is complete, this gives $\rho B^c = 0$, which implies $\rho B = 1$. \square

(P:Dynk-sp-certainty) **Proposition 4.3.5.** *Let (S, Δ, ρ) be a complete Dynkin space and let $\{A_n\} \subseteq \Delta$. If $\rho A_n = 0$ for all n , then $\rho \bigcup A_n = 0$. If $\rho A_n = 1$ for all n , then $\rho \bigcap A_n = 1$.*

Proof. Assume $\rho A_n = 0$ for all n . Let $B_n = A_n \setminus \bigcup_1^{n-1} A_j$. Then $B_n \subseteq A_n$. Since (S, Δ, ρ) is complete, we have $\rho B_n = 0$. Since $\{B_n\}$ is pairwise disjoint, we also have $\bigcup B_n \in \Delta$. But $\bigcup B_n = \bigcup A_n$, so $\rho \bigcup A_n = \rho \bigcup B_n = \sum \rho B_n = 0$. For the second claim, apply the first to A_n^c . \square

(L:ThD-is-theory) **Lemma 4.3.6.** *Let $\Omega = \mathbf{B}^{PV}$ and let $\Delta \subseteq \mathcal{B}^{PV}$ be a Dynkin system. Suppose $\mathcal{D} = (\Omega, \Delta, \rho)$ is a complete Dynkin space. Define*

$$\text{Th } \mathcal{D} = \{\varphi \in \mathcal{F} \mid \rho \varphi_\Omega = 1\}.$$

Then $\text{Th } \mathcal{D}$ is a consistent deductive theory.

Proof. Suppose $\text{Th } \mathcal{D} \vdash \varphi$. Choose countable $\Phi \subseteq \text{Th } \mathcal{D}$ such that $\vdash \bigwedge \Phi \rightarrow \varphi$. Then $\omega \models \bigwedge \Phi \rightarrow \varphi$ for all strict models ω . Hence, $(\bigwedge \Phi)_\Omega \subseteq \varphi_\Omega$. By Proposition 4.3.5, we have $\rho(\bigwedge \Phi)_\Omega = \rho \bigcap_{\theta \in \Phi} \theta_\Omega = 1$. Proposition 4.3.4 then implies $\rho \varphi_\Omega = 1$. Therefore, $\varphi \in \text{Th } \mathcal{D}$ and $\text{Th } \mathcal{D}$ is a deductive theory. Finally, $\rho \perp_\Omega = \rho \emptyset = 0$, so $\perp \notin \text{Th } \mathcal{D}$ and $\text{Th } \mathcal{D}$ is consistent. \square

(L:entire-from-Dynkin) **Lemma 4.3.7.** *Let \mathcal{D} and $\text{Th } \mathcal{D}$ be as in Lemma 4.3.6. Let $T_0 \subseteq \text{Th } \mathcal{D}$ be a deductive theory. Define $P \subseteq \mathcal{F}^{\text{IS}}$ so that $(X, \varphi, p) \in P$ if and only if there exists $\psi \in \mathcal{F}$ such that $T(X) = T_0 + \psi$ and*

$$\frac{\rho \varphi_\Omega \cap \psi_\Omega}{\rho \psi_\Omega} = p. \tag{4.3.1} \boxed{\text{cond-prob-Dynk}}$$

Then P is entire.

Proof. Note that if $A \in \Delta$ and $\rho A \triangle B = 0$, then $B \in \Delta$ and $\rho B = \rho A$. By adapting the proofs of Proposition 4.2.1, Corollary 4.2.2, and Lemma 4.2.3, we obtain the following results. If $T_0 + \psi = T_0 + \psi'$, then $\rho \psi_\Omega \triangle \psi'_\Omega = 0$. In particular, if $(X, \varphi, p) \in P$ and $T(X) = T_0 + \psi$, then (4.3.1) holds. As a consequence, if $(X, \varphi, p) \in P$ and $(X, \varphi, p') \in P$, then $p = p'$. Also, if $T(X) = T_0 + \psi$, then $(X, \varphi, 1) \in P$ if and only if $\psi_\Omega \in \Delta$ and $\rho \psi_\Omega \cap \varphi_\Omega^c = 0$.

Using these results, we may adapt the first part of the proof of Theorem 4.2.4 to show that P is entire. Note that in the proof of the addition rule, we must use the fact that if $A \cap B \in \Delta$ and two of the sets $A \cup B$, A , and B are in Δ , then all three sets are in Δ and $\rho A \cup B = \rho A + \rho B - \rho A \cap B$. \square

4.3.3 Entirety is not enough

The examples in this subsection illustrate the insufficiency of entire sets as a basis for inductive inference.

In Example 4.3.8 below, we use Dynkin spaces to construct a family of entire, strongly connected sets, indexed by $q \in (0, 1)$. Every member of this family is an example showing that probabilities of conjunctions need not be defined. That is, if P is one of the entire sets constructed in Examples 4.3.8, then $P(\mathbf{r}_1)$ exists, $P(\mathbf{r}_2)$ exists, but $P(\mathbf{r}_1 \wedge \mathbf{r}_2)$ does not exist.

We follow up in Example 4.3.10 by considering the case $q = 1/4$. In this case, we show that P is inconsistent. That is, it cannot be extended to an inductive theory. Since P is entire, it exhibits no violations whatsoever of rules (R1)–(R7). However, if we try to extend P so that it is closed under (R8), then we will inevitably create a violation of (R1)–(R7).

In Proposition 4.3.11, we consider the case $q = 1/2$. As mentioned above, $P(\mathbf{r}_1 \wedge \mathbf{r}_2)$ does not exist. Proposition 4.3.11 is concerned with what happens when we try to assign a value, q' , to this expression. The result is that $q' = 1/4$ is the unique value that we may choose in order to avoid violating rules (R1)–(R7). As such, this provides an example of the necessary use of rule (R8) to infer a probability.

(Expl:MathSE) **Example 4.3.8.** In this example, we construct an entire, strongly connected set P with root T_0 such that the domain of $P(\cdot | T_0)$ is not closed under conjunctions.

For $n \in \mathbb{N}_0$ and $k \geq 1$, define

$$d_k(n) = \lfloor 2^{-k+1}n \rfloor - 2 \lfloor 2^{-k}n \rfloor,$$

where $\lfloor x \rfloor$ is the greatest integer less than or equal to x . Then $d_k(n)$ denotes the k -th binary digit of n , counting digits from the right.

Let $PV = \{\mathbf{r}_1, \mathbf{r}_2\}$. Let $\Omega = \mathcal{B}^{PV}$ be the set of all strict models, so that $\Omega = \{\omega_n \mid 0 \leq n \leq 3\}$, where $\omega_n \mathbf{r}_k = d_k(n)$. Note that these four strict models correspond to the usual rows of a truth table with two propositional variables.

Let

$$\begin{aligned} A_1 &= (\mathbf{r}_1)_\Omega = \{\omega_1, \omega_3\}, \text{ and} \\ A_2 &= (\mathbf{r}_2)_\Omega = \{\omega_2, \omega_3\}. \end{aligned}$$

Note that

$$(\mathbf{r}_1 \leftrightarrow \mathbf{r}_2) = (A_1 \triangle A_2)^c = \{\omega_0, \omega_3\}.$$

Let $\Gamma = \{A_1, A_2, (A_1 \triangle A_2)^c\}$ and

$$\Delta = \{\emptyset, \Omega\} \cup \Gamma \cup \{A^c \mid A \in \Gamma\}.$$

Then $\Delta \subseteq \mathcal{B}^{PV}$ is a Dynkin system on Ω . Fix $q \in (0, 1)$ and, for $A \in \Gamma$, define $\rho^q A = q$ and $\rho^q A^c = 1 - q$. Together with $\rho^q \emptyset = 0$ and $\rho^q \Omega = 1$, this makes $\mathcal{D} = (\Omega, \Delta, \rho)$ a complete Dynkin space.

Let $T_0 = \mathit{Taut}$ and let P be the entire set defined in Lemma 4.3.7. Note that

$$\begin{aligned} P(\mathbf{r}_1 \mid T_0) &= q, \\ P(\mathbf{r}_2 \mid T_0) &= q, \\ P(\mathbf{r}_1 \leftrightarrow \mathbf{r}_2 \mid T_0) &= q. \end{aligned}$$

On the other hand, $(\mathbf{r}_1 \wedge \mathbf{r}_2)_\Omega = \{\omega_3\} \notin \Delta$ so that $P(\mathbf{r}_1 \wedge \mathbf{r}_2 \mid T_0)$ is undefined. Hence, P is an entire set such that the domain of $P(\cdot \mid T_0)$ is not closed under conjunctions.

$\langle \text{L:MathSE} \rangle$ **Lemma 4.3.9.** *Let $P \subseteq \mathcal{F}^{\text{IS}}$ be entire. Let $X \in \text{ante } P$ and let $\mathcal{G} \subseteq \mathcal{F}$ denote the domain of $P(\cdot \mid X)$. Let $\mathbf{r}_1, \mathbf{r}_2 \in PV$ and assume $\mathbf{r}_1, \mathbf{r}_2, \mathbf{r}_1 \leftrightarrow \mathbf{r}_2 \in \mathcal{G}$. Let*

$$\begin{aligned} \varphi_0 &= \neg \mathbf{r}_1 \wedge \neg \mathbf{r}_2, \\ \varphi_1 &= \mathbf{r}_1 \wedge \neg \mathbf{r}_2, \\ \varphi_2 &= \neg \mathbf{r}_1 \wedge \mathbf{r}_2, \\ \varphi_3 &= \mathbf{r}_1 \wedge \mathbf{r}_2. \end{aligned}$$

If $\varphi_3 \in \mathcal{G}$, then $\varphi_j \in \mathcal{G}$ for all j .

Proof. Note that $\varphi_3 \vdash \mathbf{r}_1, \mathbf{r}_2, \mathbf{r}_1 \leftrightarrow \mathbf{r}_2$. Also, $\mathbf{r}_1 \wedge \neg \varphi_3 \equiv \varphi_1$, $\mathbf{r}_2 \wedge \neg \varphi_3 \equiv \varphi_2$, and $(\mathbf{r}_1 \leftrightarrow \mathbf{r}_2) \wedge \neg \varphi_3 \equiv \varphi_0$. Hence, the result follows from Proposition 3.2.5 and the rule of logical equivalence. \square

$\langle \text{Expl:MathSEexpl-2} \rangle$ **Example 4.3.10.** Let P be as in Example 4.3.8 with $q = 1/4$. Then P is entire but not consistent. More specifically, P cannot be extended to a deductive theory. To see this, suppose P' is an inductive theory with $P \subseteq P'$. Let $P'_0 = P' \upharpoonright_{T_0}$ so that $P \subseteq P'_0$ and P'_0 is a pre-theory. Since P'_0 is a pre-theory, it is semi-closed and therefore has a completion. Let \overline{P}_0 be a completion of P'_0 . Then \overline{P}_0 is also a completion of P . Let $\mathcal{G} \subseteq \mathcal{F}$ be the domain of $\overline{P}_0(\cdot \mid T_0)$. Definition 3.3.1 implies that \mathcal{G} is closed under conjunctions. Thus, $\mathbf{r}_1 \wedge \mathbf{r}_2 \in \mathcal{G}$, so that by Lemma 4.3.9, we have $\mathbf{r}_1 \wedge \neg \mathbf{r}_2, \neg \mathbf{r}_1 \wedge \mathbf{r}_2, \neg \mathbf{r}_1 \wedge \neg \mathbf{r}_2 \in \mathcal{G}$. From Lemma 3.2.17 and Proposition 3.2.5, it follows that

$$\begin{aligned} 1 &= \overline{P}_0(\mathbf{r}_1 \mid T_0) + \overline{P}_0(\neg \mathbf{r}_1 \wedge \mathbf{r}_2 \mid T_0) + \overline{P}_0(\neg \mathbf{r}_1 \wedge \neg \mathbf{r}_2 \mid T_0) \\ &\leq \overline{P}_0(\mathbf{r}_1 \mid T_0) + \overline{P}_0(\mathbf{r}_2 \mid T_0) + \overline{P}_0(\mathbf{r}_1 \leftrightarrow \mathbf{r}_2 \mid T_0) \\ &= 3/4, \end{aligned}$$

a contradiction.

Note that this contradiction does not depend on \overline{P}_0 being complete. It is enough to assume that \overline{P}_0 is an entire extension of P such that $\overline{P}_0(\mathbf{r}_1 \wedge \mathbf{r}_2 \mid T_0)$ exists. Consequently, there is no value of q' such that $P \cup \{(T_0, \mathbf{r}_1 \wedge \mathbf{r}_2, q')\}$ has an entire extension.

$\langle \text{P:MathSE} \rangle$ **Proposition 4.3.11.** *Let P be as in Example 4.3.8 with $q = 1/2$. Let $q' \in [0, 1]$ and $Q = P \cup \{(T_0, \mathbf{r}_1 \wedge \mathbf{r}_2, q')\}$. Then Q has an entire extension if and only if $q' = 1/4$.*

Proof. Suppose P' is an entire set with $Q \subseteq P'$. Let \mathcal{G} denote the domain of $P'(\cdot | T_0)$. Let φ_j be defined as in Lemma 4.3.9, so that $\varphi_j \in \mathcal{G}$ for all j . Let $p_j = P'(\varphi_j | T_0)$ and note that $p_3 = q'$. By the addition rule,

$$p_0 + p_3 = p_1 + p_3 = p_2 + p_3 = 1/2,$$

which implies $p_0 + p_1 + p_2 + 3p_3 = 3/2$. But $\sum_j p_j = 1$, so $1 + 2p_3 = 3/2$, giving $q' = p_3 = 1/4$.

Now assume $q' = 1/4$. Let $\mathcal{P} = (\Omega, \Sigma, \mathbb{P})$, where $\Omega = \{\omega_n \mid 0 \leq n \leq 3\}$ as in Example 4.3.8, $\Sigma = \mathfrak{P}\Omega$, and \mathbb{P} satisfies $\mathbb{P}\omega_n = 1/4$ for all n . Then $\mathbf{Th} \mathcal{P}$ is an entire extension of Q . \square

We conclude this section with an example of an inductive theory that fails to be complete by violating Definition 3.3.1(ii). That is, we construct an inductive theory P where $X \in \text{ante } P$ and $X \cup \{\varphi\} \in \text{ante } P$ even though $P(\varphi | X)$ does not exist.

(Expl:incompl-ind-th)

Example 4.3.12. Let $PV = \{\mathbf{r}_1, \mathbf{r}_2, \mathbf{r}_3\}$. Recall the notation $d_k(n)$ defined in Example 4.3.8. Let Ω be the set of all strict models, so that $\Omega = \{\omega_n : 0 \leq n \leq 7\}$, where $\omega_n \mathbf{r}_k = d_k(n)$. Let $\Sigma = \mathfrak{P}\Omega$. Fix $q \in (0, 1)$ and let \mathbb{P}^q be the probability measure on (Ω, Σ) determined by

$$\begin{aligned} \mathbb{P}^q \omega_4 &= \mathbb{P}^q \omega_5 = 0, \\ \mathbb{P}^q \omega_6 &= \mathbb{P}^q \omega_7 = q/2, \\ \mathbb{P}^q \omega_n &= (1 - q)/4 \text{ for } 0 \leq n \leq 3. \end{aligned}$$

Define the model $\mathcal{P}^q = (\Omega, \Sigma, \mathbb{P}^q)$ and let $P^q = \mathbf{Th} \mathcal{P}^q$. Then P^q is a complete inductive theory with root \mathbf{Taut} . Note that

$$\begin{aligned} (\mathbf{r}_1)_\Omega &= \{\omega_1, \omega_3, \omega_5, \omega_7\}, \\ (\mathbf{r}_2)_\Omega &= \{\omega_2, \omega_3, \omega_6, \omega_7\}, \\ (\mathbf{r}_3)_\Omega &= \{\omega_4, \omega_5, \omega_6, \omega_7\}. \end{aligned}$$

Hence, $P^q(\mathbf{r}_1) = 1/2$, $P^q(\mathbf{r}_2 | \mathbf{r}_3) = 1$, and $P^q(\mathbf{r}_3) = q$.

Let Q be defined by $Q(\mathbf{r}_1) = 1/2$ and $Q(\mathbf{r}_2 | \mathbf{r}_3) = 1$. Then Q is strongly connected with root \mathbf{Taut} . Also, $Q \subseteq P^q$, so Q is consistent. Let $P = \mathbf{P}_Q$ be the inductive theory generated by Q . Then $P \subseteq P^q$ for all $q \in (0, 1)$. Hence, $P(\mathbf{r}_3)$ is undefined, and P violates Definition 3.3.1(ii).

4.4 Counterexamples and resolutions II

(S:Examples2)

This section contains examples related to satisfiability and the consequence relation. In Section 4.4.1, we construct an example where $\mathcal{P} \models \mathcal{C}$ implies $\mathcal{P} \models (X, \varphi, p)$ for all models \mathcal{P} , but $\mathcal{C} \not\models (X, \varphi, p)$. The failure occurs because X is so large that it is not countably axiomatizable over $[\mathbf{Taut}, T_{\mathcal{C}}]$. As such, this example demonstrates the need for Definition 4.2.20(ii). It also serves as

an example of a collection of inductive theories whose intersection is not an inductive theory, as well as an example of an indeterminate inductive condition.

In Section 4.4.2, we address the issue of completeness and strict satisfiability. In [16], Karp showed that completeness fails when we try to use strict satisfiability as a basis for our semantics. She presented therein two examples. In Examples 4.4.5 and 4.4.6, we revisit Karp's examples, demonstrate their resolution in the current context, and show how they connect to classical probability theory.

4.4.1 An unknown false statement

(S:D-conseq-need2) Let $PV = \{\mathbf{r}^t \mid t \in [0, 1]\}$. The idea behind this example is the following. We wish to build an inductive theory based on the following assumptions. With probability $1/2$, every propositional variable is true. Otherwise, there is exactly one $r \in PV$ that is false. But in this latter case, we do not want to make any assumptions about which r is false.

Let $T_0 = \mathbf{Taut}$. For $t \in [0, 1]$, let

$$Q(t) = \{(T_0, \mathbf{r}^t, 1/2)\} \cup \{(T_0, \mathbf{r}^s, 1) \mid s \neq t\},$$

(L:D-conseq-need2) **Lemma 4.4.1.** *For each $t \in [0, 1]$, the set $Q(t)$ is consistent.*

Proof. The set $Q(t)$ is clearly strongly connected with root T_0 . By Theorem 4.2.7, it suffices to show that $Q(t)$ is satisfiable.

Let $\Omega = \mathbf{B}^{PV}$ be the set of all strict models and $\Sigma = \mathcal{B}^{PV} = \{\varphi_\Omega \mid \varphi \in \mathcal{F}\}$. For $A \subseteq [0, 1]$, define the strict model ω_A by $\omega_A \mathbf{r}^t = 0$ if and only if $t \in A$. Let $\omega_t = \omega_{\{t\}}$. Let $\delta(w)$ be the point mass measure concentrated on w . Define $\mathbb{P}_t = (\delta(\omega_t) + \delta(\omega_\emptyset))/2$ and $\mathcal{S}_t = (\Omega, \Sigma, \mathbb{P}_t)$. Then $\omega_\emptyset \in r_\Omega$ for all $r \in PV$, and $\omega_t \in r_\Omega$ if and only if $r = \mathbf{r}^s$ for some $s \neq t$. Hence, $\bar{\mathbb{P}}_t \mathbf{r}^t = 1/2$ and $\bar{\mathbb{P}}_t \mathbf{r}^s = 1$ for $s \neq t$. Therefore, $\mathcal{S}_t \models Q(t)$ and $Q(t)$ is satisfiable. \square

Define the consistent inductive condition $\mathcal{C} = \{\mathbf{P}_{Q(t)} \mid t \in [0, 1]\}$, so that $\mathbf{P}_\mathcal{C}$ is the inductive theory we were aiming to build.

(P:D-conseq-need2) **Proposition 4.4.2.** *With notation as above, $\mathbf{P}_\mathcal{C} \subset \bigcap \mathcal{C}$. In particular, $\mathbf{P}_\mathcal{C} \notin \mathcal{C}$. That is, the condition \mathcal{C} is indeterminate.*

Proof. By Theorem 3.5.7, we have $\mathbf{P}_\mathcal{C} \subseteq \bigcap \mathcal{C}$. Assume $\mathbf{P}_\mathcal{C} = \bigcap \mathcal{C}$.

Note that for any inductive theory P , by the rule of logical implication, $X \in \text{ante } P$ if and only if $(X, \top, 1) \in P$. Hence, $\text{ante } \mathbf{P}_\mathcal{C} = \bigcap \{\text{ante } \mathbf{P}_{Q(t)} \mid t \in [0, 1]\}$.

For $t \in [0, 1]$, let $S_t = \{\mathbf{r}^s \mid s \neq t\}$. Then $\mathbf{P}_{Q(t)}(r \mid T_0) = 1$ for all $r \in S_t$ and $\mathbf{P}_{Q(t)}(\mathbf{r}^t \mid T_0) = 1/2$. By Lemma 3.2.10, we have $\{\mathbf{r}^t\} \in \text{ante } \mathbf{P}_{Q(t)}$. By the rule of deductive extension, $S_t \cup \{\mathbf{r}^t\} = PV \in \text{ante } \mathbf{P}_{Q(t)}$. Since t was arbitrary, this gives $PV \in \text{ante } \mathbf{P}_\mathcal{C}$. Hence, we may write $PV \equiv T + \psi$, where $T \in [\mathbf{Taut}, T_\mathcal{C}]$ and $\psi \in \mathcal{F}$.

Let $\Omega, \Sigma, \omega_t, \omega_\emptyset$, and \mathcal{S}_t be as in the proof of Lemma 4.4.1. Let f be the Boolean function that ψ represent. Proposition 4.1.5 implies that f is Π -ary, where $\Pi = PV \cap \text{Sf } \psi$ is the countable set of propositional variables that appear

in ψ . Hence, we may choose a measurable $h : \mathbf{B}^\Pi \rightarrow \mathbf{B}$ such that, for all $\omega \in \Omega$, we have $\omega\psi = f\omega = h(\omega|_\Pi)$.

Enumerate Π as $\Pi = \{\mathbf{r}^{t_1}, \mathbf{r}^{t_2}, \dots\}$. Choose $t_0 \notin \{t_1, t_2, \dots\}$ and let $\mathbb{P} = \delta(\omega_{t_0})$ and $\mathcal{S} = (\Omega, \Sigma, \mathbb{P})$.

Assume for the moment that $\mathcal{S} \models T_C + \psi$. Since $T + \psi \subseteq T_C + \psi$ and $PV \equiv T + \psi$, this gives $\mathcal{S} \models PV$. In particular, $\mathcal{S} \models \mathbf{r}^{t_0}$, so that $\mathbf{r}_\Omega^{t_0} \in \bar{\Sigma}$ and $\bar{\mathbb{P}}\mathbf{r}_\Omega^{t_0} = 1$. By the definition of \mathbb{P} , this implies $\omega_{t_0} \in \mathbf{r}_\Omega^{t_0}$, so that $\omega_{t_0}\mathbf{r}^{t_0} = 1$, a contradiction. Therefore, $\mathbf{P}_C \subset \bigcap \mathcal{C}$, and we are done.

It suffices, then, to show that $\mathcal{S} \models T_C + \psi$. We first show that $\mathcal{S} \models T_C$. Let $\theta \in T_C$ be arbitrary. By Proposition 3.5.8, we have $\theta \in T(\mathbf{P}_{Q(t)})$ for all $t \in [0, 1]$. Hence, $\bar{\mathbb{P}}_t \theta_\Omega = 1$ for all $t \in [0, 1]$. This implies that $\omega_\emptyset \in \theta_\Omega$ and $\omega_t \in \theta_\Omega$ for all t . In other words, $\omega_\emptyset \theta = 1$ and $\omega_t \theta = 1$ for all t .

Let $\Pi' = PV \cap \text{Sf } \theta = \{\mathbf{r}^{s_1}, \mathbf{r}^{s_2}, \dots\}$. As above, we may choose measurable $g : \mathbf{B}^{\Pi'} \rightarrow \mathbf{B}$ such that, for all $\omega \in \Omega$, we have $\omega\theta = g(\omega|_{\Pi'})$. Let $\theta^0 = \bigwedge_{j=1}^\infty \mathbf{r}^{s_j}$. For $n \in \mathbb{N}$, let $\theta^n = \neg \mathbf{r}^{s_n} \wedge (\bigwedge_{j \neq n} \mathbf{r}^{s_j})$. Finally, define $\theta' = \bigvee_{n=0}^\infty \theta^n$. If $t_0 \notin \{s_1, s_2, \dots\}$, then $\omega_{t_0}\mathbf{r}^{s_j} = 1$ for all $j \in \mathbb{N}$, which implies $\omega_{t_0}\theta^0 = 1$. If $t_0 = s_n$ for some $n \in \mathbb{N}$, then $\omega_{t_0}\theta^n = 1$. In either case, we have $\omega_{t_0}\theta' = 1$, so that $\omega_{t_0} \in \theta'_\Omega$, and therefore, $\bar{\mathbb{P}}\theta'_\Omega = 1$.

Now suppose $\omega \in \theta'_\Omega$. Choose $n \in \mathbb{N}_0$ such that $\omega \in \theta^n_\Omega$. If $n = 0$, then $\omega|_{\Pi'} = \omega_\emptyset|_{\Pi'}$, so that $\omega\theta = \omega_\emptyset\theta = 1$ and $\omega \in \theta_\Omega$. If $n \in \mathbb{N}$, then $\omega|_{\Pi'} = \omega_{s_n}|_{\Pi'}$, so that $\omega\theta = \omega_{s_n}\theta = 1$ and again $\omega \in \theta_\Omega$. This shows that $\theta'_\Omega \subseteq \theta_\Omega$. Therefore, $\bar{\mathbb{P}}\theta_\Omega = 1$, so that $\mathcal{S} \models \theta$. Since θ was arbitrary, we have $\mathcal{S} \models T_C$.

Lastly, we show that $\mathcal{S} \models \psi$. Since $(\Omega, \Sigma, \delta(\omega_\emptyset)) \models PV$ and $PV \equiv T + \psi$, we have $\omega_\emptyset \in \psi_\Omega$, so that $\omega_\emptyset\psi = 1$. Since $t_0 \notin \{t_1, t_2, \dots\}$, we also have that $\omega_\emptyset|_\Pi = \omega_{t_0}|_\Pi$. Hence, $\omega_{t_0}\psi = \omega_\emptyset\psi = 1$, which gives $\omega_{t_0} \in \psi_\Omega$ and therefore $\mathcal{S} \models \psi$. \square

(R:intersect-th-fail)

Remark 4.4.3. Since \mathbf{P}_C is the largest inductive theory contained in $\bigcap \mathcal{C}$, it follows that $\bigcap \mathcal{C}$ is not an inductive theory. By Theorem 3.5.3, the set $\bigcap \mathcal{C}$ is closed. Hence, it must not be connected. In other words, the condition \mathcal{C} is an example of a collection of connected sets whose intersection is not connected.

(R:D-conseq-need2)

Remark 4.4.4. Note that if $\mathcal{S} \models \mathcal{C}$, then $\mathcal{S} \models \mathbf{P}_{Q(t)}$ for some t . But we also have $(PV, \top, 1) \in \mathbf{P}_{Q(t)}$ for all t . Hence, $\mathcal{S} \models \mathcal{C}$ implies $\mathcal{S} \models (PV, \top, 1)$. On the other hand, the proof of Proposition 4.4.2 shows that $PV \notin \text{ante } \mathbf{P}_C$. Therefore, this example illustrates the necessity of Definition 4.2.20(ii).

4.4.2 Karp's counterexamples

(Σξάρξαρξ13)

Example 4.4.5. It is well-known that σ -compactness fails for strict satisfiability. That is, there exists $X \subseteq \mathcal{F}$ such that every countable subset is strictly satisfiable, but X itself is not strictly satisfiable. Karp gives an example of such an X in [16, Example 4.1.3].

However, since strict satisfiability implies satisfiability (Proposition 4.1.6(i)) and satisfiability is σ -compact (Theorem 4.1.17), we know that any such X

must be satisfiable. We present Karp's example below, and then show that it is satisfied by one of the most common models in probability theory.

Let $PV = \{\mathbf{r}_n^k \mid (n, k) \in \mathbb{N} \times \{0, 1\}\}$. Let $\zeta = \bigwedge_n (\mathbf{r}_n^0 \vee \mathbf{r}_n^1)$, and for each $f : \mathbb{N} \rightarrow \{0, 1\}$, let $\psi_f = \neg \bigwedge_n \mathbf{r}_n^{f(n)}$. Let $X = \{\zeta\} \cup \{\psi_f \mid f \in \{0, 1\}^{\mathbb{N}}\}$.

Suppose $X_0 \subseteq X$ is countable. Since $\{0, 1\}^{\mathbb{N}}$ is uncountable, there exists $g : \mathbb{N} \rightarrow \{0, 1\}$ such that $\psi_g \notin X_0$. Define ω by $\omega \mathbf{r}_n^k = 1$ if and only if $g(n) = k$. Then $\omega \models \zeta$. Note that $\omega \models \psi_f$ if and only if there exists n such that $\omega \mathbf{r}_n^{f(n)} = 0$. Given $f \neq g$, we may choose n such that $f(n) \neq g(n)$, and so for this value of n , we have $\omega \mathbf{r}_n^{f(n)} = 0$. Hence, $\omega \models \psi_f$ for all $f \neq g$, and therefore $\omega \models X_0$. It follows that every countable subset of X is strictly satisfiable.

Now suppose $\omega \models X$. Since $\omega \models \zeta$, we may choose, for each $n \in \mathbb{N}$, a value $f(n) \in \{0, 1\}$ such that $\omega \models \mathbf{r}_n^{f(n)}$. But then $\omega \models \bigwedge_n \mathbf{r}_n^{f(n)}$, meaning $\omega \not\models \psi_f$, a contradiction. Therefore, X is not strictly satisfiable.

As mentioned in the beginning of this example, however, we know that there exists a model $\mathcal{P} = (\Omega, \Sigma, \mathbb{P})$ such that $\mathcal{P} \models X$. In this case, we can construct such a model using $\Omega = \mathcal{B}^{PV}$ and taking $\Sigma = \mathcal{B}^{PV} = \{\varphi_\Omega \mid \varphi \in \mathcal{F}\}$, thereby assigning a probability to every formula in \mathcal{F} . The model is a natural one that is ubiquitous in probability theory. Namely, it is the one that models an i.i.d. sequence of coin flips.

Let (S, Γ, ν) be a probability space on which we have constructed an i.i.d. sequence $\langle X_n \mid n \in \mathbb{N} \rangle$ of $\{0, 1\}$ -valued random variables with $\nu\{X_n = 1\} = 1/2$. Define $G : PV \rightarrow \Gamma$ by $G\mathbf{r}_n^k = \{X_n = k\}$. Let $\mathcal{P} = (\Omega, \Sigma, \mathbb{P})$ be the model constructed in the proof of Theorem 4.3.1. Note that

$$\begin{aligned} G\zeta &= \{X_n \in \{0, 1\} \text{ for all } n\}, \\ G\psi_f &= \{X_n = f(n) \text{ for all } n\}^c. \end{aligned}$$

Hence, $\mathbb{P}\zeta_\Omega = \nu G\zeta = 1$ and $\mathbb{P}(\psi_f)_\Omega = \nu G\psi_f = 1$ for all f , showing that $\mathcal{P} \models X$.

We will investigate this example further in Section 4.5.5, after covering the topic of independence.

(Exp1:Karp412) **Example 4.4.6.** We present here another example of Karp's (see [16, Example 4.1.2]). Again we see an X that demonstrates the failure of σ -compactness for strict satisfiability. And again, we know that X is satisfiable. We could use the construction in the proof of Theorem 4.1.17 to build a model that satisfies X . In this example, though, we do not do that. Rather, we show that any such model has a certain property. Namely, in any such model, $\mathcal{P} = (\Omega, \Sigma, \mathbb{P})$, there will be formulas $\varphi \in \mathcal{F}$ such that $\varphi_\Omega \notin \Sigma$.

Let I be an uncountable set and let $PV = \{\mathbf{r}_n^t \mid t \in I, n \in \mathbb{N}\}$. For each $t \in I$, let $\zeta^t = \bigvee_n \mathbf{r}_n^t$. For each $s, t \in I$ and $n \in \mathbb{N}$, let $\psi_n^{s,t} = \neg(\mathbf{r}_n^s \wedge \mathbf{r}_n^t)$. Then define

$$X = \{\zeta^t \mid t \in I\} \cup \{\psi_n^{s,t} \mid s, t \in I, s \neq t, n \in \mathbb{N}\}.$$

Let $X_0 \subseteq X$ be countable. Then there is a countable set $S \subset I$ such that

$$X_0 \subseteq \{\zeta^t \mid t \in S\} \cup \{\psi_n^{s,t} \mid s, t \in S, s \neq t, n \in \mathbb{N}\}.$$

Let $t \mapsto n(t)$ be an injection from S to \mathbb{N} . Define a model ω by $\omega \mathbf{r}_n^t = 1$ if and only if $t \in S$ and $n = n(t)$. Then, for each $t \in S$, we have $\omega \models \mathbf{r}_{n(t)}^t$, and therefore $\omega \models \zeta^t$. Also, if $s, t \in S$, $s \neq t$, and $n \in \mathbb{N}$, then either $n \neq n(s)$ or $n \neq n(t)$, implying that at least one of $\omega \mathbf{r}_n^s$ and $\omega \mathbf{r}_n^t$ is 0. Thus, $\omega \models \psi_n^{s,t}$. Altogether, this shows that $\omega \models X_0$, so that every countable subset of X is strictly satisfiable.

Now suppose $\omega \models X$. For each $t \in I$, we have $\omega \models \zeta^t$. By the definition of ζ^t , we may choose $n(t) \in \mathbb{N}$ such that $\omega \models \mathbf{r}_{n(t)}^t$. But I is uncountable and \mathbb{N} is countable, so there exists distinct $s, t \in I$ such that $n(s) = n(t)$. It follows that $\omega \psi_{n(s)}^{s,t} = 0$, and so $\omega \not\models \psi_{n(s)}^{s,t}$, a contradiction. Therefore, X is not strictly satisfiable.

Again, by Proposition 4.1.6(i) and Theorem 4.1.17, we know that X is satisfiable, meaning there is a model $\mathcal{P} = (\Omega, \Sigma, \mathbb{P})$ such that $\mathcal{P} \models X$. In this case, however, we cannot construct such a model using $\Sigma = \mathcal{B}^{PV} = \{\varphi_\Omega \mid \varphi \in \mathcal{F}\}$. In any model that satisfies X , there will exist $\varphi \in \mathcal{F}$ such that $\varphi_\Omega \notin \Sigma$. In other words, there will be formulas that are not assigned a probability.

To see that this is the case, let $\mathcal{P} = (\Omega, \Sigma, \mathbb{P})$ be a model. Assume that $\mathcal{P} \models X$ and that $r_\Omega \in \Sigma$ for all $r \in PV$. For $n, k \in \mathbb{N}$, define

$$S(n, k) = \{t \in I \mid \mathbb{P}(\mathbf{r}_n^t)_\Omega \geq k^{-1}\}.$$

Suppose $s, t \in I$ and $s \neq t$. Then $\psi_n^{s,t} = \neg(\mathbf{r}_n^s \wedge \mathbf{r}_n^t) \in X$. Thus,

$$\mathbb{P}(\psi_n^{s,t})_\Omega = \mathbb{P}((\mathbf{r}_n^s)_\Omega \cap (\mathbf{r}_n^t)_\Omega)^c = 1.$$

In other words, $(\mathbf{r}_n^s)_\Omega$ and $(\mathbf{r}_n^t)_\Omega$ are pairwise disjoint, up to a set of measure zero. It follows that $S(n, k)$ is a finite set with at most k elements.

Now fix $t \in I$. Since $\zeta^t = \bigvee_n \mathbf{r}_n^t \in X$, we have $1 = \mathbb{P} \zeta^t = \mathbb{P} \bigcup_n (\mathbf{r}_n^t)_\Omega$. Thus, there exists $n \in \mathbb{N}$ such that $\mathbb{P}(\mathbf{r}_n^t)_\Omega > 0$, showing that $t \in S(n, k)$ for some $n, k \in \mathbb{N}$. In other words, $I = \bigcup_{n,k} S(n, k)$, expressing I as a countable union of finite sets, contradicting the fact that I is uncountable.

^(R:Karp412) **Remark 4.4.7.** It follows from Example 4.4.6 that we cannot construct an \mathbb{N} -valued stochastic process $\langle Y(t) \mid t \in I \rangle$ such that for all $s \neq t$, we have $Y(s) \neq Y(t)$ a.s. (Recall that a stochastic process is simply an indexed collection of random variables taking values in the same measurable space.) To see this, suppose we have such a process, built on a probability space (S, Γ, ν) . Define $G : PV \rightarrow \Gamma$ by $G \mathbf{r}_n^t = \{Y(t) = n\}$, and let $\mathcal{P} = (\Omega, \Sigma, \mathbb{P})$ be the model constructed in the proof of Theorem 4.3.1. Then $\mathcal{P} \models X$ and $\Sigma = \{\varphi_\Omega \mid \varphi \in \mathcal{F}\}$, a contradiction.

4.5 Independence

^(S:indep) In this section, we introduce the concept of (inductive) independence. It is a purely logical concept, defined solely in terms of formulas and inductive theories, without reference to any model. Intuitively, φ and ψ are independent given X

if $P(\varphi \mid X)$ is unchanged by adding ψ to X . More generally, a sequence of formulas is independent if, whenever we take two disjoint subsequences and, from each subsequence, create a formula using negation and conjunction, the two resulting formulas are independent.

To make this notion precise, we will introduce the concept of a dialog set, which describes all the formulas that can be created from a given set of formulas. Dialog sets are the syntactic analogues of σ -algebras.

After defining independence, we introduce the semantic concept of “measure independence,” which is nothing more than the familiar notion of independence in a measure space, defined by the property that the measure of an intersection factors into a product. Using this, we give a characterization of independence in terms of satisfiability and measure independence.

Finally, we give two examples in which we use independence to build inductive conditions, and then look at the inductive theories generated by those conditions.

4.5.1 Dialog sets

(B:dialog) **Definition 4.5.1.** A set $D \subseteq \mathcal{F}$ is a *dialog set* if it satisfies the following:

- (i) $\varphi \in D$ implies $\neg\varphi \in D$,
- (neg-cl)** (ii) $\Phi \subseteq D$ countable implies $\bigwedge \Phi \in D$, and
- (conj-cl)** (iii) $\varphi \in D$ and $\varphi \equiv \psi$ implies $\psi \in D$.
- (eq-cl)**

Note that the Φ in (ii) may be empty. Hence, $\top \in D$, and by (i) and (iii), also $\perp \in D$. Note further that (iii) implies D is closed under \rightarrow , \leftrightarrow , and \bigvee .

The intersection of any family of dialog sets is again a dialog set. Also, the language \mathcal{F} itself is a dialog set, and is the largest dialog set. If $X \subseteq \mathcal{F}$, then the *dialog set generated by X* , denoted by $\delta(X)$, is the smallest dialog set containing X . It is equal to the intersection of all dialog sets containing X . The smallest dialog set is

$$\delta(\emptyset) = \text{Taut} \cup \neg\text{Taut}.$$

Intuitively, $\delta(X)$ is the set of all formulas that can be built out the formulas in X using negation, conjunction, and logical equivalence. Note that $\delta(PV) = \mathcal{F}$.

(P:dialog-sig-alg) **Proposition 4.5.2.** Let Ω be a set of strict models. If $D \subseteq \mathcal{F}$ is a dialog set, then D_Ω is a σ -algebra on Ω . More generally, if $X \subseteq \mathcal{F}$, then $\sigma(X_\Omega) = \delta(X)_\Omega$.

Proof. Let $D \subseteq \mathcal{F}$ be a dialog set. Then $\Omega = \top_\Omega \in D_\Omega$. Let $A \in D_\Omega$. Choose $\varphi \in D$ such that $A = \varphi_\Omega$. Then $\neg\varphi \in D$, so that $A^c = \varphi_\Omega^c = (\neg\varphi)_\Omega \in D_\Omega$. Finally, suppose $\{A_n\} \subseteq D_\Omega$. Choose $\varphi_n \in D$ such that $A_n = (\varphi_n)_\Omega$. Then $\bigcap A_n = (\bigwedge \varphi_n)_\Omega \in D_\Omega$, and hence, D_Ω is a σ -algebra.

Now let $X \subseteq \mathcal{F}$. Then $X_\Omega \subseteq \delta(X)_\Omega$. By the above, $\delta(X)_\Omega$ is a σ -algebra. Hence, $\sigma(X_\Omega) \subseteq \delta(X)_\Omega$. For the reverse inclusion, define $D = \{\varphi \in \mathcal{F} \mid \varphi_\Omega \in \sigma(X_\Omega)\}$. The set D clearly satisfies (i) and (ii) of Definition 4.5.1. Remark 4.1.14 shows that it also satisfies (iii). Thus, D is a dialog set. Since $X \subseteq D$, it follows that $\delta(X) \subseteq D$ and therefore $\delta(X)_\Omega \subseteq \sigma(X_\Omega)$. \square

4.5.2 Independence of two formulas

Let P be an inductive theory. For $X \subseteq \mathcal{F}$, let $\text{dom } P(\cdot | X) \subseteq \mathcal{F}$ denote the domain of $P(\cdot | X)$. If $X \notin \text{ante } P$, then $\text{dom } P(\cdot | X) = \emptyset$.

Let $\varphi, \psi \in \text{dom } P(\cdot | X)$. We say that φ is *dependent on ψ given X (under P)* if $P(\varphi | X, \psi)$ exists and is not equal to $P(\varphi | X)$. We say that φ is *independent of ψ given X (under P)* if either $P(\psi | X) = 0$ or $P(\varphi | X, \psi) = P(\varphi | X)$.

Note that if $P(\psi | X) > 0$ and $P(\varphi | X, \psi)$ does not exist, then φ is not dependent on ψ , and φ is also not independent of ψ . If P is complete, then this situation cannot arise. To see this, note that if P is complete and $\varphi, \psi \in \text{dom } P(\cdot | X)$, then $P(\varphi \wedge \psi | X)$ exists. Hence, by the multiplication rule, if $P(\psi | X) > 0$, then $P(\varphi | X, \psi)$ exists.

$\langle \text{L:dep-undec} \rangle$ **Lemma 4.5.3.** *If φ is dependent on ψ given X , then both $0 < P(\varphi | X) < 1$ and $0 < P(\psi | X) < 1$.*

Proof. Let φ be dependent on ψ given X . By Lemma 3.2.10, we have $P(\psi | X) > 0$. Suppose $P(\psi | X) = 1$. Then the multiplication rule and Proposition 3.2.11 imply $P(\varphi | X) = P(\varphi | X, \psi)$, a contradiction.

Now suppose $P(\varphi | X) = 1$. As above, the multiplication rule and Proposition 3.2.11 imply $P(\varphi | X, \psi) = 1$, a contradiction. Finally, suppose $P(\varphi | X) = 0$. Then $P(\neg\varphi | X) = 1$. Again, this gives $P(\neg\varphi | X, \psi) = 1$, so that (3.2.5) implies $P(\varphi | X, \psi) = 0$, a contradiction. \square

$\langle \text{L:evidence} \rangle$ **Lemma 4.5.4.** *If φ is dependent on ψ given X , then both $P(\psi | X, \varphi)$ and $P(\varphi | X, \neg\psi)$ exist.*

Proof. Suppose φ is dependent on ψ given X . Then $P(\varphi | X)$, $P(\psi | X)$, and $P(\varphi | X, \psi)$ exist, and $P(\varphi | X, \psi) \neq P(\varphi | X)$. Since $P(\psi | X)$ and $P(\varphi | X, \psi)$ exist, the multiplication rule implies $P(\varphi \wedge \psi | X)$ exists. By Lemma 4.5.3, we have $P(\varphi | X) > 0$. Hence, another application of the multiplication rule implies that $P(\psi | X, \varphi)$ exists. From Proposition 3.2.5, it follows that $P(\varphi \wedge \neg\psi | X)$ exists. Lemma 4.5.3 implies $P(\neg\psi | X) > 0$. Therefore, a final application of the multiplication rule gives the existence of $P(\varphi | X, \neg\psi)$. \square

Proposition 4.5.5. *Suppose φ is dependent on ψ given X . Then $P(\varphi | X, \psi) > P(\varphi | X)$ if and only if $P(\varphi | X, \neg\psi) < P(\varphi | X)$.*

Proof. Suppose $P(\varphi | X, \psi) > P(\varphi | X)$. Lemma 4.5.3 implies $P(\psi | X) > 0$, and Lemma 4.5.4 implies $P(\psi | X, \varphi)$ exists. Thus, by (3.2.8),

$$P(\varphi | X)P(\psi | X, \varphi) > P(\psi | X)P(\varphi | X),$$

which implies $P(\psi | X, \varphi) > P(\psi | X)$. By (3.2.5), this implies $P(\neg\psi | X, \varphi) < P(\neg\psi | X)$. On the other hand, Lemma 4.5.4 implies $P(\varphi | X, \neg\psi)$ exists, so that (3.2.8) implies

$$P(\neg\psi | X)P(\varphi | X, \neg\psi) = P(\varphi | X)P(\neg\psi | X, \varphi).$$

As before, this implies $P(\varphi \mid X, \neg\psi) < P(\varphi \mid X)$. The analogous argument proves the converse. \square

$\langle \text{P:depen-sym} \rangle$ **Proposition 4.5.6.** *Let $\varphi, \psi \in \mathcal{F}$. Then φ is dependent on ψ given X if and only if ψ is dependent on φ given X .*

Proof. Let φ be dependent on ψ given X . Then $P(\varphi \mid X)$ and $P(\psi \mid X)$ exist. By Lemma 4.5.4, we have that $P(\psi \mid X, \varphi)$ exists. Suppose $P(\psi \mid X, \varphi) = P(\psi \mid X)$. Lemma 4.5.3 implies $P(\psi \mid X) > 0$. Thus, by (3.2.8), we have $P(\varphi \mid X) = P(\varphi \mid X, \psi)$, a contradiction. Therefore, ψ is dependent on φ given X . The converse follows by reversing the roles of φ and ψ . \square

$\langle \text{T:indep-prod} \rangle$ **Theorem 4.5.7.** *Let $\varphi, \psi \in \text{dom } P(\cdot \mid X)$. Then φ is independent of ψ given X if and only if*

$$P(\varphi \wedge \psi \mid X) = P(\varphi \mid X)P(\psi \mid X). \quad (4.5.1) \quad \boxed{\text{indep-prod}}$$

Proof. Suppose $P(\psi \mid X) = 0$. Then φ is independent of ψ given X . Also, Proposition 3.2.11 implies $P(\varphi \wedge \psi \mid X) = 0$, so that (4.5.1) holds.

Now suppose $P(\psi \mid X) > 0$. Then φ is independent of ψ given X if and only if $P(\varphi \mid X, \psi) = P(\varphi \mid X)$. On the other hand, by the multiplication rule, (4.5.1) holds if and only if $P(\varphi \mid X, \psi) = P(\varphi \mid X)$. \square

$\langle \text{C:indep-prod} \rangle$ **Corollary 4.5.8.** *Let $\varphi, \psi \in \text{dom } P(\cdot \mid X)$. If φ is independent of ψ given X , then φ is independent of $\neg\psi$ given X .*

Proof. Suppose φ is independent of ψ given X . By Proposition 3.2.5 and Theorem 4.5.7, we have

$$\begin{aligned} P(\varphi \wedge \neg\psi \mid X) &= P(\varphi \mid X) - P(\varphi \wedge \psi \mid X) \\ &= P(\varphi \mid X) - P(\varphi \mid X)P(\psi \mid X) \\ &= P(\varphi \mid X)(1 - P(\psi \mid X)) \\ &= P(\varphi \mid X)P(\neg\psi \mid X). \end{aligned}$$

Hence, Theorem 4.5.7 implies φ is independent of $\neg\psi$ given X . \square

By Proposition 4.5.6 and Theorem 4.5.7, we can alter our terminology to say that φ and ψ are dependent or independent, given X . Note that by Theorem 4.5.7 and Proposition 3.2.11, if $\varphi, \psi \in \text{dom } P(\cdot \mid X)$ and either $P(\varphi \mid X) \in \{0, 1\}$ or $P(\psi \mid X) \in \{0, 1\}$, then φ and ψ are independent given X .

4.5.3 Independence of a sequence of formulas

Let I be a set with $|I| \geq 2$ and let $\langle \varphi_i \mid i \in I \rangle$ be an indexed collection of formulas in $\text{dom } P(\cdot \mid X)$. Such a collection is *independent given X (under P)* if φ and ψ are independent given X whenever $\varphi \in \delta(\{\varphi_i \mid i \in I_1\})$ and $\psi \in \delta(\{\varphi_i \mid i \in I_2\})$, where I_1 and I_2 are nonempty disjoint subsets of I .

Proposition 4.5.9. *Let $\varphi, \psi \in \text{dom } P(\cdot \mid X)$. Then φ and ψ are independent given X if and only if $\langle \varphi, \psi \rangle$ is independent given X .*

Proof. The if direction is trivial. For the only if direction, suppose φ and ψ are independent given X and let $\varphi' \in \delta(\{\varphi\})$ and $\psi' \in \delta(\{\psi\})$. Note that $\delta(\{\varphi\})$ consists of tautologies, contradictions, and formulas that are equivalent to either φ or $\neg\varphi$, and similarly for $\delta(\{\psi\})$. We may assume that $0 < P(\varphi' | X) < 1$ and $0 < P(\psi' | X) < 1$, so that neither φ nor ψ is a tautology or contradiction.

Clearly, φ' and ψ' are independent if $\varphi' \equiv \varphi$ and $\psi' \equiv \psi$. By Corollary 4.5.8, φ' and ψ' are independent if $\varphi' \equiv \varphi$ and $\psi' \equiv \neg\psi$. Repeated applications of this result cover the cases $\varphi' \equiv \neg\varphi$ and $\psi' \equiv \neg\psi$, and $\varphi' \equiv \neg\varphi$ and $\psi' \equiv \psi$. \square

$\langle \text{T:indep-dialog-defined} \rangle$ **Theorem 4.5.10.** *Let P be an inductive theory and $X \in \text{ante } P$. If $\langle \varphi_i | i \in I \rangle$ is independent given X , then $\delta(\{\varphi_i | i \in I\}) \subseteq \text{dom } P(\cdot | X)$.*

Proof. Let $\Omega = \mathbf{B}^{PV}$ and define Δ as in (4.1.2), so that $\Delta = Y_\Omega$, where $Y = \text{dom } P(\cdot | X)$. Proposition 4.1.16 implies that Δ is a Dynkin system on Ω . Let $U = \{(\varphi_i)_\Omega | i \in I\}$. Since $\langle \varphi_i | i \in I \rangle$ is independent given X , it follows that $U \subseteq \Delta$. By Theorem 4.5.7, we have that U is a π -system, that is, U is closed under pairwise intersections. Therefore, Dynkin's π - λ theorem gives $\sigma(U) \subseteq \Delta$.

Now let $\varphi \in \delta(\{\varphi_i | i \in I\})$. By Proposition 4.5.2, we have $\varphi_\Omega \in \sigma(U) \subseteq \Delta$. Hence, we may choose $\varphi' \in \text{dom } P(\cdot | X)$ such that $\varphi_\Omega = \varphi'_\Omega$. Since $\Omega = \mathbf{B}^{PV}$, according to Remark 4.1.14, it follows that $\varphi \equiv \varphi'$. Therefore, by the rule of logical equivalence, $\varphi \in \text{dom } P(\cdot | X)$. \square

4.5.4 A semantic characterization of independence

$\langle \text{S:sem-indep} \rangle$ Let (S, Γ, ν) be a probability space and I a set with $|I| \geq 2$. For each $i \in I$, let $A_i \in \Gamma$. Then $\langle A_i | i \in I \rangle$ is *measure independent* in (S, Γ, ν) if $\nu \bigcap_{i \in J} A_i = \prod_{i \in J} \nu A_i$, whenever $J \subseteq I$ is finite. Note that this is the usual definition of independence in a probability space. Also note that we may assume without loss of generality that $|J| \geq 2$.

Let P be an inductive theory. Let $\mathcal{P} = (\Omega, \Sigma, \mathbb{P})$ be a model and suppose $\mathcal{P} \models P$. Let $X \in \text{ante } P$. Write $X \equiv Y \cup \{\psi\}$, where $\mathcal{P} \models Y$ and $\mathbb{P}\psi_\Omega > 0$. Define the probability measure $\overline{\mathbb{P}}_X$ on $(\Omega, \overline{\Sigma})$ by $\overline{\mathbb{P}}_X A = \overline{\mathbb{P}} A \cap \psi_\Omega / \overline{\mathbb{P}}\psi_\Omega$, and let $\mathcal{P}_X = (\Omega, \overline{\Sigma}, \overline{\mathbb{P}}_X)$. Note that by Proposition 4.2.1, the model \mathcal{P}_X does not depend on our choice of Y and ψ .

$\langle \text{T:indep-sound-compl} \rangle$ **Theorem 4.5.11.** *Let P be an inductive theory, $X \in \text{ante } P$, and I a set with $|I| \geq 2$. Let $\langle \varphi_i | i \in I \rangle$ be an indexed collection of formulas in $\text{dom } P(\cdot | X)$. Then the following are equivalent:*

- (i) $\langle \varphi_i | i \in I \rangle$ is independent given X .
- (ii) For any model \mathcal{P} , if $\mathcal{P} \models P$, then $\langle (\varphi_i)_\Omega | i \in I \rangle$ is measure independent in \mathcal{P}_X .
- (iii) There exists a model \mathcal{P} such that $\mathcal{P} \models P$ and $\langle (\varphi_i)_\Omega | i \in I \rangle$ is measure independent in \mathcal{P}_X .

Proof. Suppose (i) holds. Assume \mathcal{P} is a model and $\mathcal{P} \models P$. Let $J \subseteq I$ be finite with $|J| \geq 2$. Fix $k \in J$. Let $I_1 = \{k\}$ and $I_2 = J \setminus \{k\}$. Then $\varphi_k \in \delta(\{\varphi_i : i \in I_1\})$ and $\bigwedge_{i \in I_2} \varphi_i \in \delta(\{\varphi_i : i \in I_2\})$ are independent given X . By (4.5.1),

$$P(\bigwedge_{i \in J} \varphi_i \mid X) = P(\varphi_k \mid X)P(\bigwedge_{i \in J \setminus \{k\}} \varphi_i \mid X).$$

Since $\mathcal{P} \models P$, this implies

$$\bar{\mathbb{P}}_X \bigcap_{i \in J} (\varphi_i)_\Omega = \bar{\mathbb{P}}_X (\varphi_k)_\Omega \bar{\mathbb{P}}_X \bigcap_{i \in J \setminus \{k\}} (\varphi_i)_\Omega.$$

Iterating this argument gives $\bar{\mathbb{P}}_X \bigcap_{i \in J} (\varphi_i)_\Omega = \prod_{i \in J} \bar{\mathbb{P}}_X (\varphi_i)_\Omega$, so that the indexed collection $\langle (\varphi_i)_\Omega \mid i \in I \rangle$ is measure independent in \mathcal{P}_X , showing that (i) implies (ii).

By Theorem 4.2.6, (ii) implies (iii).

Suppose (iii) holds. For $i \in I$, let $A_i = (\varphi_i)_\Omega$, so that by hypothesis, $\langle A_i \mid i \in I \rangle$ is measure independent in \mathcal{P}_X . A result from measure theory tells us that if I_1 and I_2 are disjoint subsets of I , then A and B are measure independent in \mathcal{P}_X whenever $A \in \sigma(\{A_i \mid i \in I_1\})$ and $B \in \sigma(\{A_i \mid i \in I_2\})$.

Let I_1 and I_2 be nonempty disjoint subsets of I . Let $U = \{\varphi_i \mid i \in I_1\}$ and $V = \{\varphi_i \mid i \in I_2\}$. Let $\varphi \in \delta(U)$ and $\psi \in \delta(V)$. By Proposition 4.5.2, we have $\varphi_\Omega \in \sigma(U_\Omega)$ and $\psi_\Omega \in \sigma(V_\Omega)$, so that φ_Ω and ψ_Ω are measure independent in \mathcal{P}_X . Hence,

$$\bar{\mathbb{P}}_X \varphi_\Omega \cap \psi_\Omega = \bar{\mathbb{P}}_X \varphi_\Omega \bar{\mathbb{P}}_X \psi_\Omega. \quad (4.5.2) \quad \boxed{\text{indep-sound-compl}}$$

Theorem 4.5.10 implies φ , ψ , and $\varphi \wedge \psi$ are all in $\text{dom } P(\cdot \mid X)$. Since $\mathcal{P} \models P$, it follows that (4.5.2) implies $P(\varphi \wedge \psi \mid X) = P(\varphi \mid X)P(\psi \mid X)$. By Theorem 4.5.7, therefore, φ and ψ are independent given X . Thus, (iii) implies (i). \square

$\langle \text{C: indep-sound-compl} \rangle$ **Corollary 4.5.12.** *Let P be an inductive theory, $X \in \text{ante } P$, and I a set with $|I| \geq 2$. Let $\langle \varphi_i \mid i \in I \rangle$ be an indexed collection of formulas in $\text{dom } P(\cdot \mid X)$. Then $\langle \varphi_i \mid i \in I \rangle$ is independent given X if and only if*

$$P(\bigwedge_{j \in J} \varphi_j \mid X) = \prod_{j \in J} P(\varphi_j \mid X)$$

for all finite $J \subseteq I$.

Proof. This follows immediately from Theorems 4.5.11 and 4.2.6. \square

4.5.5 Fair coin flips

$\langle \text{S: Karp413} \rangle$ In this subsection, our aim is to create an inductive theory that describes an infinite sequence of independent flips of a fair coin.

We must first construct the language in which this will be done. Let $PV = \{\mathbf{r}_n^k \mid (n, k) \in \mathbb{N} \times \{0, 1\}\}$. We interpret \mathbf{r}_n^k as representing the proposition, “The n th flip of the coin lands on k .” Here, $k = 1$ represents heads and $k = 0$ represents tails.

Our inductive theory will be built on three “axioms,” informally stated as:

- (1) Each flip must land on heads or tails.

- (2) On an individual flip, the probabilities of heads and tails are each $1/2$.
- (3) The flips are independent.

We will enforce (1) with our choice of root, T_0 . We will enforce (2) with a set Q of inductive statements. We will enforce (3) with an inductive condition, \mathcal{C} .

Let $T_0 = T(\{\zeta\})$, where $\zeta = \bigwedge_n (\mathbf{r}_n^0 \vee \mathbf{r}_n^1)$. Let

$$Q = \{(T_0, \mathbf{r}_n^k, 1/2) \mid (n, k) \in \mathbb{N} \times \{0, 1\}\}.$$

Note that Q is connected with root T_0 . Define the inductive condition \mathcal{C} to be the set of all inductive theories with root T_0 such that $\langle \mathbf{r}_n^{f(n)} \mid n \in \mathbb{N} \rangle$ is independent given T_0 whenever $f \in \{0, 1\}^{\mathbb{N}}$.

Recall that $Q, \mathcal{C} \vdash (X, \varphi, p)$ means $\mathcal{C}_Q \cap \mathcal{C} \vdash (X, \varphi, p)$.

(L:fair-coin-flips) **Lemma 4.5.13.** *The inductive condition $\mathcal{C}_Q \cap \mathcal{C}$ is consistent.*

Proof. Let $\mathcal{P} = (\Omega, \Sigma, \mathbb{P})$ be the model constructed in Example 4.4.5. As shown in that example, $\mathcal{P} \models T_0$. Hence, by Proposition 3.5.10, we have that $P = \mathbf{Th} \mathcal{P} \downarrow_{[T_0, Th \mathcal{P}]}$ is an inductive theory with root T_0 . Note that

$$\mathbb{P}(\mathbf{r}_n^k)_\Omega = \nu Gr_n^k = \nu\{X_n = k\} = 1/2.$$

Hence, $P(\mathbf{r}_n^k \mid T_0) = \overline{\mathbb{P}}(\mathbf{r}_n^k)_\Omega = 1/2$, so that $Q \subseteq P$, and therefore $P \in \mathcal{C}_Q$.

Let $f \in \{0, 1\}^{\mathbb{N}}$. Since $(r_n^{f(n)})_\Omega = \{X_n = f(n)\}$, it follows that $\langle (r_n^{f(n)})_\Omega \mid n \in \mathbb{N} \rangle$ is measure independent in \mathcal{P} . Thus, Theorem 4.5.11 implies $P \in \mathcal{C}$, and so $P \in \mathcal{C}_Q \cap \mathcal{C}$. Hence, $\mathcal{C}_Q \cap \mathcal{C}$ is nonempty, and therefore consistent. \square

By Lemma 4.5.13, we may define $\mathbf{P}_{Q, \mathcal{C}} = \mathbf{P}(\mathcal{C}_Q \cap \mathcal{C})$, the inductive theory generated by Q and \mathcal{C} . Note that $Q, \mathcal{C} \vdash (X, \varphi, p)$ if and only if $\mathbf{P}_{Q, \mathcal{C}}(\varphi \mid X) = p$. In other words, $\mathbf{P}_{Q, \mathcal{C}}$ is precisely the inductive theory we are aiming for. It contains exactly those inductive statements that can be derived from (1)–(3).

(P:fair-coin-flips-deter) **Proposition 4.5.14.** *The inductive condition $\mathcal{C}_Q \cap \mathcal{C}$ is determinate. That is, $\mathbf{P}_{Q, \mathcal{C}} \in \mathcal{C}_Q \cap \mathcal{C}$.*

Proof. Suppose $\mathcal{P} \models \mathcal{C}_Q \cap \mathcal{C}$. Choose $P \in \mathcal{C}_Q \cap \mathcal{C}$ such that $\mathcal{P} \models P$. Then P is an inductive theory with root T_0 such that $Q \subseteq P$ and $\langle \mathbf{r}_n^{f(n)} \mid n \in \mathbb{N} \rangle$ is independent given T_0 whenever $f \in \{0, 1\}^{\mathbb{N}}$. Let $f \in \{0, 1\}^{\mathbb{N}}$ and let $I \subseteq \mathbb{N}$ be finite. By Corollary 4.5.12, we have $P(\bigwedge_{i \in I} \mathbf{r}_i^{f(i)} \mid T_0) = 2^{-|I|}$. Hence, $\mathcal{P} \models (T_0, \bigwedge_{i \in I} \mathbf{r}_i^{f(i)}, 2^{-|I|})$. By Theorem 4.2.21, we have $Q, \mathcal{C} \vdash (T_0, \bigwedge_{i \in I} \mathbf{r}_i^{f(i)}, 2^{-|I|})$. Therefore, $\mathbf{P}_{Q, \mathcal{C}}(\bigwedge_{i \in I} \mathbf{r}_i^{f(i)} \mid T_0) = 2^{-|I|}$. Again by Corollary 4.5.12, this shows that $\mathbf{P}_{Q, \mathcal{C}} \in \mathcal{C}$. Taking $|I| = 1$ shows $\mathbf{P}_{Q, \mathcal{C}} \in \mathcal{C}_Q$. \square

(P:fair-coin-flips-char) **Proposition 4.5.15.** *Let P be the inductive theory in the proof of Lemma 4.5.13. Then $P = \mathbf{P}_{Q, \mathcal{C}}$.*

Proof. Since $P \in \mathcal{C}_Q \cap \mathcal{C}$, we have $\mathbf{P}_{Q,\mathcal{C}} \subseteq \bigcap \mathcal{C}_Q \cap \mathcal{C} \subseteq P$. For the reverse inclusion, since P and $\mathbf{P}_{Q,\mathcal{C}}$ both have root T_0 , it suffices to show that $P \downarrow_{T_0} \subseteq \mathbf{P}_{Q,\mathcal{C}} \downarrow_{T_0}$. For notational simplicity, let $P' = \mathbf{P}_{Q,\mathcal{C}}$.

We first show that $P'(\varphi \mid T_0)$ exists for every $\varphi \in \mathcal{F}$. Let $Y \subseteq \mathcal{F}$ be the set of all finite conjunctions of propositional variables. We claim that $Y \subseteq \text{dom } P'(\cdot \mid T_0)$. To see this, let $\varphi \in Y$. Suppose, for some n , that φ contains both \mathbf{r}_n^0 and \mathbf{r}_n^1 . Then $\varphi \vdash \mathbf{r}_n^0 \wedge \mathbf{r}_n^1$. Since $\mathbf{r}_n^0 \vee \mathbf{r}_n^1 \in T_0$, we have $P'(\mathbf{r}_n^0 \vee \mathbf{r}_n^1 \mid T_0) = 1$. Hence, Theorem 3.2.18 implies $P'(\mathbf{r}_n^0 \wedge \mathbf{r}_n^1 \mid T_0) = 0$, and therefore, $P'(\varphi \mid T_0) = 0$. On the other hand, suppose that φ contains at most one of \mathbf{r}_n^0 and \mathbf{r}_n^1 for each n . Then $\varphi = \bigwedge_{i \in I} \mathbf{r}_i^{f(i)}$ for some finite $I \subseteq \mathbb{N}$ and some $f \in \{0, 1\}^{\mathbb{N}}$. By the proof of Proposition 4.5.14, we have $P'(\varphi \mid T_0) = 2^{-|I|}$. Hence, $Y \subseteq \text{dom } P'(\cdot \mid T_0)$.

Recall that in Lemma 4.5.13, we have $\Omega = \mathbf{B}^{PV}$. The above shows that $Y_\Omega \subseteq \Delta$, where $\Delta = \{\varphi_\Omega \mid P'(\varphi \mid T_0) \text{ exists}\}$. Proposition 4.1.16 implies that Δ is a Dynkin system. Since Y_Ω is closed under pairwise intersections, Dynkin's π - λ theorem implies that $\mathbf{B}^{PV} = \sigma(Y_\Omega) \subseteq \Delta$. Hence, if $\varphi \in \mathcal{F}$, then $\varphi_\Omega \in \Delta$, so that $\varphi_\Omega = \varphi'_\Omega$ for some $\varphi' \in \text{dom } P'(\cdot \mid T_0)$. Remark 4.1.14 gives $\varphi \equiv \varphi'$, so that by the rule of logical equivalence, $P'(\varphi \mid T_0)$ exists.

Now suppose $P(\varphi \mid T_0, \psi) = p$. Then, by the definition of P , we have $\overline{\mathbb{P}} \varphi_\Omega \cap \psi_\Omega / \overline{\mathbb{P}} \psi_\Omega = p$. Since $P'(\varphi \wedge \psi \mid T_0)$ and $P'(\psi \mid T_0)$ both exist and $\mathcal{P} \models P'$, this implies that $P'(\varphi \wedge \psi \mid T_0) / P'(\psi \mid T_0) = p$. From the multiplication rule, it follows that $P'(\varphi \mid T_0, \psi) = p$. \square

Recall $\psi_f = \neg \bigwedge_n \mathbf{r}_n^{f(n)}$ from Example 4.4.5, where $f \in \{0, 1\}^{\mathbb{N}}$. The function f is simply a sequence of 1's and 0's. If we interpret f as a sequence of heads and tails, then the formula ψ_f represents the sentence,

“The pattern of heads and tails produced by the coin is not f .”

By Proposition 4.5.15 and Example 4.4.5, we have $\mathbf{P}_{Q,\mathcal{C}}(\psi_f \mid T_0) = 1$ for every $f \in \{0, 1\}^{\mathbb{N}}$. Hence, $T_0, Q, \mathcal{C} \vdash \psi_f$, so that ψ_f is a logical consequence of (1)–(3), and this is true for every $f \in \{0, 1\}^{\mathbb{N}}$.

Classical intuition suggests that this is paradoxical, since the coin must produce *some* pattern. But this classical intuition is rooted in the idea of strict satisfiability. Indeed, Example 4.4.5 shows that there is no strict model that strictly satisfies both T_0 and every ψ_f . A strict model is an assignment of truth values to every sentence. Classical intuition thinks in terms of these truth assignments. To remove any cognitive dissonance produced by this example, intuition must be changed so that it thinks in terms of probability measures on truth assignments.

4.5.6 Biased coin flips

As in the previous subsection, our aim here is to create an inductive theory that describes an infinite sequence of independent coin flips. This time, however, we will drop the assumption that the coin is fair.

As before, we use the language built on $PV = \{\mathbf{r}_n^k \mid (n, k) \in \mathbb{N} \times \{0, 1\}\}$, where \mathbf{r}_n^k represents the proposition, “The n th flip of the coin lands on k .”

This time, the “axioms” of our inductive theory will be:

- (1) Each flip must land on heads or tails.
- (2) On an individual flip, the probabilities of heads and tails sum to 1.
- (3) Every flip has the same probability of heads, which is neither 0 nor 1.
- (4) The flips are independent.

We will enforce (1) with our choice of root, T_0 . We will enforce (2)–(4) with inductive conditions.

As before, let $T_0 = T(\{\zeta\})$, where $\zeta = \bigwedge_n (\mathbf{r}_n^0 \vee \mathbf{r}_n^1)$. Let \mathfrak{J}_{T_0} be the set of inductive theories with root T_0 . Define the inductive conditions,

$$\mathcal{C}_2 = \{P \in \mathfrak{J}_{T_0} \mid P(\mathbf{r}_n^0 \mid T_0) + P(\mathbf{r}_n^1 \mid T_0) = 1 \text{ for all } n\},$$

$$\mathcal{C}_q = \{P \in \mathfrak{J}_{T_0} \mid P(\mathbf{r}_n^1 \mid T_0) = q \text{ for all } n\}, \quad q \in (0, 1),$$

$$\mathcal{C}_4 = \{P \in \mathfrak{J}_{T_0} \mid \langle \mathbf{r}_n^{f(n)} \mid n \in \mathbb{N} \rangle \text{ is independent given } T_0 \text{ for all } f \in \{0, 1\}^{\mathbb{N}}\}$$

Let $\mathcal{C}_3 = \bigcup_{q \in (0, 1)} \mathcal{C}_q$. Then \mathcal{C}_j represents assumption (j) for $2 \leq j \leq 4$. Let $\mathcal{C} = \mathcal{C}_2 \cap \mathcal{C}_3 \cap \mathcal{C}_4$.

(P:bias-coin-flips) **Proposition 4.5.16.** *The condition \mathcal{C} is consistent, but indeterminate. That is, $\mathbf{P}_{\mathcal{C}} \notin \mathcal{C}$. More precisely, the domain of $\mathbf{P}_{\mathcal{C}}(\cdot \mid T_0)$ does not contain any propositional variables. Hence, $\mathbf{P}_{\mathcal{C}} \notin \mathcal{C}_2$, $\mathbf{P}_{\mathcal{C}} \notin \mathcal{C}_4$, and $\mathbf{P}_{\mathcal{C}} \notin \mathcal{C}_q$ for any $q \in (0, 1)$.*

Proof. Let \mathcal{P} be the model constructed in Example 4.4.5 and let $P = \mathbf{Th} \mathcal{P} \downarrow_{[T_0, Th \mathcal{P}]}$. Then $P \in \mathcal{C}_2 \cap \mathcal{C}_{1/2} \cap \mathcal{C}_4 \subseteq \mathcal{C}$ and $\mathcal{P} \models P$. Hence, $\mathcal{P} \models \mathcal{C}$, so that \mathcal{C} is satisfiable and therefore consistent.

Fix $r \in PV$ and assume $\mathbf{P}_{\mathcal{C}}(r \mid T_0)$ exists. Let $q_0 = \mathbf{P}_{\mathcal{C}}(r \mid T_0)$. Choose $q \in (0, 1)$ such that $q_0 \notin \{q, 1 - q\}$. As in Example 4.4.5, we may construct a model \mathcal{P}^q such that $P^q = \mathbf{Th} \mathcal{P}^q \downarrow_{[T_0, Th \mathcal{P}^q]} \in \mathcal{C}_2 \cap \mathcal{C}_q \cap \mathcal{C}_4 \subseteq \mathcal{C}$. Since $P^q \in \mathcal{C}_q$, we have $P^q(r \mid T_0) \in \{q, 1 - q\}$. On the other hand, since $P^q \in \mathcal{C}$, it follows that $\mathbf{P}_{\mathcal{C}} \subseteq \bigcap \mathcal{C} \subseteq P^q$, so that $P^q(r \mid T_0) = q_0$, a contradiction. Hence, $\mathbf{P}_{\mathcal{C}}(r \mid T_0)$ does not exist. \square

Proposition 4.5.16 shows that the domain of $\mathbf{P}_{\mathcal{C}}(\cdot \mid T_0)$ does not contain any propositional variables. This domain, however, is not trivial. That is, it contains more than just tautologies and contradictions. Recall the formulas $\psi_f = \neg \bigwedge_n \mathbf{r}_n^{f(n)}$, where $f \in \{0, 1\}^{\mathbb{N}}$. As in Example 4.4.5, we can show that $P(\psi_f \mid T_0) = 1$ for every $P \in \mathcal{C}$. Since $\mathbf{P}_{\mathcal{C}} \downarrow_{T_0} = \bigcap \mathcal{C}^0$, it follows that $\mathbf{P}_{\mathcal{C}}(\psi_f \mid T_0) = 1$ for every $f \in \{0, 1\}^{\mathbb{N}}$.

It might be tempting to think that $\mathbf{P}_{\mathcal{C}}(\cdot \mid T_0)$ is entirely deductive, in the sense that $\mathbf{P}_{\mathcal{C}}(\varphi \mid T_0) \in \{0, 1\}$ whenever $\mathbf{P}_{\mathcal{C}}(\varphi \mid T_0)$ exists. After all, the inductive condition \mathcal{C} does not specify any numerical probabilities at all. What

we will show, however, is that the opposite is true. For any $p \in (0, 1)$, there is a formula $\varphi \in \mathcal{F}$ such that $\mathbf{P}_{\mathcal{C}}(\varphi \mid T_0) = p$.

The intuition behind this is the following. It is possible to simulate a fair coin flip with a biased coin. Simply flip the coin twice. If the results match, start over. If they do not match, use the second of the two flips as the result. But we can do this as many times as we like. So we can simulate an i.i.d. sequence of fair coin flips. We can then use this sequence to simulate a random number that is uniformly chosen from the interval $(0, 1)$. Finally, we construct a formula which asserts that this uniform random number is less than p .

Proposition 4.5.17. *For any $p \in (0, 1)$, there exists a formula $\varphi \in \mathcal{F}$ such that $\mathbf{P}_{\mathcal{C}}(\varphi \mid T_0) = p$.*

Proof. Fix $p \in (0, 1)$. Let $\Omega = \mathbf{B}^{PV}$ be the set of all strict models and $\Sigma = \mathcal{B}^{PV} = \{\varphi_{\Omega} \mid \varphi \in \mathcal{F}\}$. Let $\zeta' = \bigwedge_n \neg(\mathbf{r}_n^0 \wedge \mathbf{r}_n^1)$. Define $Y_n : \Omega \rightarrow \{0, 1\}$ as follows. If $\omega \not\models \zeta \wedge \zeta'$, then $Y_n(\omega) = 0$. If $\omega \models \zeta \wedge \zeta'$, then define $Y_n(\omega)$ so that $\omega \models \mathbf{r}_n^{Y_n(\omega)}$. Note that $\{Y_n = 1\} = (\zeta \wedge \zeta' \wedge \mathbf{r}_n^1)_{\Omega} \in \Sigma$. Hence, Y_n is Σ -measurable.

Define $\tau_0 = 0$ and

$$\tau_k = \inf\{n > \tau_{k-1} \mid n \text{ is even and } Y_n \neq Y_{n-1}\}.$$

Define $Z_k = 0$ on $\{\tau_k = \infty\}$ and $Z_k = Y_{\tau_k}$ on $\{\tau_k < \infty\}$, and let $U = \sum_1^{\infty} 2^{-k} Z_k$. Then U is Σ -measurable, which implies $\{U \leq p\} \in \Sigma$. Choose $\varphi \in \mathcal{F}$ such that $\{U \leq p\} = \varphi_{\Omega}$. We will show that $\mathbf{P}_{\mathcal{C}}(\varphi \mid T_0) = p$.

Fix $P \in \mathcal{C}$. Let \mathcal{P} be an arbitrary model with $\mathcal{P} \models P$. Define $P' = \mathbf{Th} \mathcal{P} \downarrow_{[T_0, Th \mathcal{P}]}$, so that P' is a complete inductive theory with root T_0 such that $P \subseteq P'$. By the proof of Theorem 4.2.6, we may construct a model $\mathcal{P}' = (\Omega, \Sigma, \mathbb{P})$, where $\Omega = \mathbf{B}^{PV}$ and $\Sigma = \mathcal{B}^{PV}$, such that $P' = \mathbf{Th} \mathcal{P}' \downarrow_{[T_0, Th \mathcal{P}]}$.

Using $P \in \mathcal{C}$, $P \subseteq P'$, and Theorem 3.2.18, we have $P'(\zeta' \mid T_0) = 1$. Hence, $\{Y_n = 1\} = (\mathbf{r}_n^1)_{\Omega}$ \mathbb{P} -a.s. Thus, $\{Y_n\}$ are i.i.d. in \mathcal{P}' with $\mathbb{P}\{Y_n = 1\} = q$, where q is the value satisfying $P \in \mathcal{C}_q$. It is a straightforward exercise to verify that $\{Z_k\}$ are i.i.d. with $\mathbb{P}\{Z_k = 1\} = 1/2$. Hence, U is uniformly distributed on $(0, 1)$, which gives $P'(\varphi \mid T_0) = \mathbb{P} \varphi_{\Omega} = \mathbb{P}\{U \leq p\} = p$.

From the definition of P' , we have $\mathcal{P}' \models (T_0, \varphi, p)$. Since \mathcal{P} was arbitrary, Theorem 4.2.14 implies $P \vdash (T_0, \varphi, p)$. But P is an inductive theory, so this gives $P(\varphi \mid T_0) = p$. Finally, since P was arbitrary, it follows that $(T_0, \varphi, p) \in \bigcap \mathcal{C}^0 = \mathbf{P}_{\mathcal{C}} \downarrow_{T_0}$, so that $\mathbf{P}_{\mathcal{C}}(\varphi \mid T_0) = p$. \square

Chapter 5

Predicate Logic

(Ch:pred-logic) In this chapter, we repeat the work we did in Chapters 3 and 4, but in the setting of a predicate language. Most of the work in this chapter is devoted to the deductive side of predicate logic. The development of inductive logic requires hardly any modification from the propositional case.

The language we use is just like first-order logic, except it allows countable conjunctions and disjunctions. It is typically denoted in the literature by $\mathcal{L}_{\omega_1, \omega}$. We will denote it simply by \mathcal{L} .

In Section 5.1, we introduce the syntax of \mathcal{L} . We define terms, formulas, sentences, and all of their related concepts, such as subformulas, free and bound variables, and substitutions. The set of formulas is denoted by \mathcal{L} . As in the propositional case, formulas are built from an alphabet of symbols. To the propositional alphabet, we add an uncountable number of variable symbols, and we add the logical symbol \forall . Unlike the propositional case, our alphabet will not include the set PV . Instead, it will include a set of symbols L , called the extralogical signature. The set L includes constant symbols, relation symbols, and function symbols.

Sentences are formulas that have no free variables. For example, $x > 0$ is a formula, whereas $\forall x x > 0$ is a sentence. The set of formulas is denoted by \mathcal{L} , and the set of sentences is denoted by $\mathcal{L}^0 \subseteq \mathcal{L}$. Intuitively, a sentence says something. It can be meaningful and have a truth value. On the other hand, a formula is ambiguous. It could mean many different things, depending on the values assigned to its free variables. As such, it cannot have its own truth value. Predicate logic is concerned with sentences. In fact, both deductive and inductive theories consist exclusively of sentences. As such, our models and our inferential calculi should deal directly with \mathcal{L}^0 . In the inductive case, that is exactly what we will do. In the deductive case, however, we do something different. In that case, it will be easier for us to build models and inferential calculi for \mathcal{L} . When we do so, we will treat free variables as if they are constant symbols.

Section 5.2 is concerned with inferential calculus in predicate logic. The bulk of this section is devoted to a system of natural deduction for deductive

inference in \mathcal{L} . We present this system, prove that it is σ -compact, and then show how it connects to inference in \mathcal{L}^0 . We discuss the inductive calculus, which carries over almost entirely unchanged from Chapter 3. Finally, we give a Hilbert-type calculus for deductive inference. Specifically, this is the calculus used by Carol Karp in [16].

In Section 5.3, we present the semantics of predicate logic. As before, the majority of the section is given to deductive semantics. Inductive semantics carry over from Chapter 4 with very little modifications. In predicate logic, sentences are given meaning by interpreting them in a structure. A model will be a probability measure on a set of structures. Using this, we will define satisfiability and the consequence relation, then prove σ -compactness, soundness, and completeness.

Section 5.3 also contains the theory of Peano arithmetic in the infinitary setting. We present the usual theory from first-order logic, as well as two different extensions to the infinitary language \mathcal{L}^0 . That is, we define three theories, $\text{PA}_{\text{fin}} \subseteq \text{PA}_- \subseteq \text{PA}$. The theory PA_{fin} is the usual theory of Peano arithmetic in first-order logic. The theories PA_- and PA are extensions to \mathcal{L}^0 . The first extension, PA_- , is conservative, in the sense that every sentence in $\text{PA}_- \setminus \text{PA}_{\text{fin}}$ is purely infinitary. In other words, in PA_- , we cannot deduce any new first-order sentences that we could not already deduce in PA_{fin} . This is because PA_- and PA_{fin} have the same axioms. In particular, even in PA_- , we can only do induction on finitary formulas. The second extension, PA , is stronger. There, we allow induction on infinitary formulas. In doing so, we find that PA completely characterizes the natural numbers. That is, every true sentence about arithmetic is provable in PA .

Finally, in Section 5.4, we connect inductive logic to the measure-theoretic concept of a random variable. As described in Section 1.9, we will be forced to deal with an issue we call “the relativity of randomness.” We will do so by introducing and discussing “frames of references.” The connection between inductive logic and random variables will allow us to show that a measure-theoretic probability model—that is, a probability space, together with a collection of random variables—is a special case of a model in inductive logic. In other words, measure-theoretic probability is embedded in inductive logic. In fact, this embedding is proper. As we will see in Example 5.4.8, inductive logic is capable of expressing things that cannot be expressed in measure-theoretic probability.

5.1 The syntax of predicate formulas

(S:syntax) In this section, we present the predicate language \mathcal{L} . We describe the alphabet, and the rules for constructing terms, formulas, and sentences. In Karp’s original construction of \mathcal{L} (see [16]), formulas are countably long strings of symbols. Our construction follows [18] instead, building up formulas out of sets.

5.1.1 The alphabet and terms

$\langle \mathbf{S} : \text{terms} \rangle$ Let L be an extralogical signature. Recall the convention that, unless otherwise stated, c denotes a constant symbol, r a relation symbol, and f a function symbol with arity $n \geq 1$.

Let $\mathit{Var} = \{\mathbf{x}_\alpha \mid \alpha < \omega_1\}$ be an uncountable set of symbols. The symbols in Var are called *individual variables*. Unless otherwise stated, letters such as u, v, x, y, z will denote distinct individual variables.

We define an alphabet, $\mathbf{A} = L \cup \mathit{Var} \cup \{\neg, \wedge, \forall, =\}$. Let $\mathcal{S} = \mathcal{S}_L$ denote the set of (finite) strings over \mathbf{A} . We use boldface for the symbol, $=$, to distinguish it from the ordinary equal sign. For instance, $\xi = \mathbf{x}_0 = \mathbf{x}_1$ means that ξ is equal to the length-3 string, $\mathbf{x}_0 = \mathbf{x}_1$. Parentheses are not part of our alphabet, but we may sometimes add them for readability. For example, the above might be written as $\xi = (\mathbf{x}_0 = \mathbf{x}_1)$ to further emphasize the distinction between $=$ and $=$. We may also write $\xi : \mathbf{x}_0 = \mathbf{x}_1$, as another way to improve readability.

Definition 5.1.1. The set of *terms* in L , denoted by $\mathcal{T} = \mathcal{T}_L$, is the smallest subset of \mathcal{S} such that

- (i) $\mathit{Var} \subseteq \mathcal{T}$,
- (ii) $c \in \mathcal{T}$ for all constant symbols $c \in L$, and
- (iii) if $f \in L$ is an n -ary function symbol and $t_1, \dots, t_n \in \mathcal{T}$, then $ft_1 \cdots t_n \in \mathcal{T}$.

Unless otherwise stated, letters such as s and t will denote terms. Individual variables and constant symbols are called *prime terms*. A term is called *compound* if it is not prime. Terms of the form in (iii) are called *function terms*. Note that every compound term is a function term.

Also note that terms have the same definition here as they do in first-order logic. Hence, they have all the same properties. For instance (see [28, Section 2.2]), the *unique term concatenation property* says that if $t_1 \cdots t_n = s_1 \cdots s_m$, then $n = m$ and $t_i = s_i$ for all i . Also, the *unique term reconstruction property* says that if $ft_1 \cdots t_n = fs_1 \cdots s_n$, then $t_i = s_i$ for all i .

When convenient, we will adopt the notation $\vec{t} = t_1 \cdots t_n$ for a concatenation of terms. We also adopt shorthand to improve the readability of terms. For example, suppose $x, y, z \in \mathit{Var}$, and let $+$ and \circ be binary operation symbols. Then $t = +xy$ is a term, and $otz = o+xyz$ is a term. This latter term is especially difficult to read, and we would typically write it as $(x + y) \circ z$. Note that parentheses are not symbols in our alphabet; this is simply shorthand. We also adopt this shorthand for relations, so that $<xy$ would be written as $x < y$.

Definition 5.1.2. For $t \in \mathcal{T}$, the set $\text{var } t \subseteq \mathit{Var}$ is defined recursively as follows:

- (i) if $c \in L$ is a constant symbol, then $\text{var } c = \emptyset$,
- (ii) if $x \in \mathit{Var}$, then $\text{var } x = \{x\}$, and
- (iii) $\text{var } ft_1 \cdots t_n = \text{var } t_1 \cup \cdots \cup \text{var } t_n$.

Intuitively, $\text{var } t$ is the set of individual variables occurring in t . By induction on t , it can be shown that $\text{var } t$ is countable for all $t \in \mathcal{T}$. If $\text{var } t = \emptyset$, then we call t a *ground term*, or *constant term*.

Definition 5.1.3. For $t \in \mathcal{T}$, the set $\text{sym } t \subseteq L$ is defined recursively as follows:

- (i) if $c \in L$ is a constant symbol, then $\text{sym } c = \{c\}$,
- (ii) if $x \in \text{Var}$, then $\text{sym } x = \emptyset$, and
- (iii) $\text{sym } ft_1 \cdots t_n = \{f\} \cup \text{sym } t_1 \cup \cdots \cup \text{sym } t_n$.

Intuitively, $\text{sym } t$ is the set of extralogical symbols occurring in t . Note that $\text{sym } t$ is also countable. In addition, we define

$$\text{con } t = \{c \in \text{sym } t \mid c \text{ is a constant symbol}\},$$

which denotes the (countable) set of constant symbols occurring in t .

5.1.2 Formulas

(S:pred-formulas)

A string $\varphi \in \mathcal{S}$ is an *equation* if $\varphi = (s = t)$, where $s, t \in \mathcal{T}$. A string $\varphi \in \mathcal{S}$ is a *prime (or atomic) formula* if it is an equation or if it has the form $\varphi = rt_1 \cdots t_n$, where $r \in L$ is an n -ary relation symbol and $t_1, \dots, t_n \in \mathcal{T}$.

Note that prime formulas have the same definition here as they do in first-order logic. Hence, they have all the same properties. For instance (see [28, Section 2.2]), the *unique prime formula reconstruction property* says that if $rt_1 \cdots t_n = rs_1 \cdots s_n$, then $t_i = s_i$ for all i . Also, terms do not contain the symbol $=$. Therefore, if $(s = t) = (s' = t')$, then $s = s'$ and $t = t'$.

We will define the set of formulas so that a formula is a finite tuple, where each element in the tuple is either a symbol from our alphabet, a formula, or a countable set of formulas.

Let S_0 denote the set of prime formulas. For an ordinal $\alpha < \omega_1$, let

$$S'_\alpha = S_\alpha \cup \{\langle \neg, \varphi \rangle \mid \varphi \in S_\alpha\} \cup \{\langle \forall, x, \varphi \rangle \mid x \in \text{Var}, \varphi \in S_\alpha\}.$$

As with strings, when writing tuples such as these, we will typically omit the commas and angled brackets, so that, for instance, $\forall x \varphi = \langle \forall, x, \varphi \rangle$.

We then define

$$S_{\alpha+1} = S'_\alpha \cup \{\langle \bigwedge, \Phi \rangle \mid \Phi \subseteq S'_\alpha \text{ is nonempty and countable}\}.$$

Here, countable means finite or countably infinite. As above, we will typically write $\bigwedge \Phi$ as a shorthand for ordered pairs of this type.

In the case that α is a nonzero limit ordinal, we define $S_\alpha = \bigcup_{\xi < \alpha} S_\xi$. Finally, we define $\mathcal{L} = \mathcal{L}_{\omega_1, \omega} = \bigcup_{\alpha < \omega_1} S_\alpha$. Note that $S_\alpha \subseteq S_\beta$ whenever $\alpha < \beta$. An element $\varphi \in \mathcal{L}$ is called a *formula*. A formula φ is called a *literal* if $\varphi = \pi$ or $\varphi = \neg \pi$ for some prime formula π .

Theorem 5.1.4 (Unique formula reconstruction property). *If φ is a formula that is not prime, then exactly one of the following holds.*

- (i) *There exists a unique $\psi \in \mathcal{L}$ such that $\varphi = \neg\psi$.*
- (ii) *There exists a unique $x \in \mathbf{Var}$ and a unique $\psi \in \mathcal{L}$ such that $\varphi = \forall x\psi$.*
- (iii) *There exists a unique $\Phi \subseteq \mathcal{L}$ such that $\varphi = \bigwedge \Phi$.*

Proof. Let $\varphi \in \mathcal{L}$ and assume φ is not prime. Let β be the smallest ordinal such that $\varphi \in S_\beta$. Since φ is not prime, $\beta > 0$. Since $\varphi \notin S_\xi$ for all $\xi < \beta$, it follows that β is not a limit ordinal. Therefore, β is a successor ordinal, and we may write $\beta = \alpha + 1$. Since $\varphi \notin S_\alpha$, we have

$$\begin{aligned} \varphi \in \{ \langle \neg, \psi \rangle \mid \psi \in S_\alpha \} \cup \{ \langle \forall, x, \psi \rangle \mid x \in \mathbf{Var}, \psi \in S_\alpha \} \\ \cup \{ \langle \bigwedge, \Phi \rangle \mid \Phi \subseteq S'_\alpha \text{ is nonempty and countable} \}. \end{aligned}$$

Note that the above union is a disjoint union. Hence, φ is in exactly one of the above three sets. \square

Let \mathcal{L}_{fin} denote the smallest subset of \mathcal{L} that satisfies

- (i) prime formulas are in \mathcal{L}_{fin} ,
- (ii) if $\varphi \in \mathcal{L}_{\text{fin}}$ and $x \in \mathbf{Var}$, then $\neg\varphi \in \mathcal{L}_{\text{fin}}$ and $\forall x\varphi \in \mathcal{L}_{\text{fin}}$, and
- (iii) if $\Phi \subseteq \mathcal{L}_{\text{fin}}$ is nonempty and finite, then $\bigwedge \Phi \in \mathcal{L}_{\text{fin}}$.

Formulas in \mathcal{L}_{fin} are said to be *finitary*. The set \mathcal{L}_{fin} is, in fact, the set of formulas used in first-order logic. The reader can consult any introductory text on mathematical logic for the basic properties of \mathcal{L}_{fin} and its corresponding syntax and semantics. When necessary, we will cite [28] for this purpose.

We adopt all the same shorthand as in the propositional language \mathcal{F} , except for the definitions of falsum and verum, which will be given later in Section 5.2.5. In addition, we also use the shorthand $\exists x\varphi = \neg\forall\neg\varphi$ and $(s \neq t) = \neg(s = t)$. We may also write $\forall x_1x_2 \cdots x_n$ or $\forall \vec{x}$ instead of $\forall x_1\forall x_2 \cdots \forall x_n$. If \triangleright is a binary relation symbol, we will write

$$\begin{aligned} (\forall x \triangleright t)\varphi &= \forall x(x \triangleright t \rightarrow \varphi), \\ (\exists x \triangleright t)\varphi &= \exists x(x \triangleright t \wedge \varphi), \end{aligned}$$

and similarly for ∇ . In Section 5.1.5, after introducing substitutions, we will give shorthand for $\exists!$, the unique existential quantifier.

The set of formulas \mathcal{L} depends on the extralogical signature L . We may sometimes emphasize this fact in our notation. For example, if $L = \{\circ, e\}$, we may write $\mathcal{L} = \mathcal{L}\{\circ, e\}$. We may also write $\mathcal{S}_{\mathcal{L}}$ and $\mathcal{T}_{\mathcal{L}}$ instead of \mathcal{S}_L and \mathcal{T}_L .

5.1.3 Formula induction and recursion

The proof of Theorem 3.1.1 carries over with minor modification to give us the following.

Theorem 5.1.5 (The principle of formula induction). *The set of formulas, \mathcal{L} , is the smallest set that satisfies the following:*

- (i) *prime formulas are in \mathcal{L} ,*
- (ii) *if $\varphi \in \mathcal{L}$ and $x \in \mathit{Var}$, then $\neg\varphi \in \mathcal{L}$ and $\forall x\varphi \in \mathcal{L}$, and*
- (iii) *if $\Phi \subseteq \mathcal{L}$ is nonempty and countable, then $\bigwedge \Phi \in \mathcal{L}$.*

Given $\varphi \in \mathcal{L}$, we define $\mathit{Sf} \varphi \subseteq \mathcal{L}$, the set of *subformulas* of φ , by formula recursion. Namely, $\mathit{Sf} \pi = \{\pi\}$ if π is prime, $\mathit{Sf} \neg\varphi = \{\neg\varphi\} \cup \mathit{Sf} \varphi$, $\mathit{Sf} \bigwedge \Phi = \{\bigwedge \Phi\} \cup \bigcup_{\varphi \in \Phi} \mathit{Sf} \varphi$, and $\mathit{Sf} \forall x\varphi = \{\forall x\varphi\} \cup \mathit{Sf} \varphi$. It follows by formula induction that $\mathit{Sf} \varphi$ is countable for every $\varphi \in \mathcal{L}$.

Given $\varphi \in \mathcal{L}$, we define $\mathit{len} \varphi \in \mathbb{N} \cup \{\infty\}$, which we call the *length* of φ , by formula recursion. If φ is prime, then φ is a finite, nonempty string of symbols from our alphabet A . In this case, let $\mathit{len} \varphi$ be the length of this string. We then extend this by $\mathit{len} \neg\varphi = 1 + \mathit{len} \varphi$, $\mathit{len} \forall x\varphi = 2 + \mathit{len} \varphi$, and $\mathit{len} \bigwedge \Phi = 1 + \sum_{\varphi \in \Phi} \mathit{len} \varphi$. Note that if φ has a subformula of infinite length, then φ has infinite length. Also note that $\varphi \in \mathcal{L}_{\text{fin}}$ if and only if $\mathit{len} \varphi < \infty$.

Given $\varphi \in \mathcal{L}$, we define the ordinal $\mathit{rk} \varphi$, called the *rank* of φ , by formula recursion. Namely, $\mathit{rk} \pi = 0$ if π is prime, $\mathit{rk} \neg\varphi = \mathit{rk} \varphi + 1$, $\mathit{rk} \bigwedge \Phi = (\bigcup_{\varphi \in \Phi} \mathit{rk} \varphi) + 1$, and $\mathit{rk} \forall x\varphi = \mathit{rk} \varphi + 1$. Note that $\mathit{rk} \varphi = 0$ if and only if φ is prime, and $\mathit{rk} \varphi$ is a successor ordinal whenever φ is not prime.

5.1.4 Variables and symbols

Definition 5.1.6. For $\varphi \in \mathcal{L}$, the set $\mathit{var} \varphi \subseteq \mathit{Var}$ is defined recursively as follows:

- (i) $\mathit{var} s = t = \mathit{var} s \cup \mathit{var} t$,
- (ii) $\mathit{var} r t_1 \cdots t_n = \mathit{var} t_1 \cup \cdots \cup \mathit{var} t_n$,
- (iii) $\mathit{var} \neg\varphi = \mathit{var} \varphi$,
- (iv) $\mathit{var} \bigwedge \Phi = \bigcup_{\varphi \in \Phi} \mathit{var} \varphi$, and
- (v) $\mathit{var} \forall x\varphi = \mathit{var} \varphi \cup \{x\}$.

Intuitively, $\mathit{var} \varphi$ is the set of individual variables occurring in φ . It follows by formula induction that $\mathit{var} \varphi$ is countable for every $\varphi \in \mathcal{L}$. In other words, even though Var is uncountable, any given formula will only make use of countably many individual variables.

Given $\varphi \in \mathcal{L}$, we define $\mathit{bnd} \varphi \subseteq \mathit{Var}$, the set of *bound variables* in φ , by formula recursion. Namely, $\mathit{bnd} \pi = \emptyset$ if π is prime, $\mathit{bnd} \neg\varphi = \mathit{bnd} \varphi$,

$\text{bnd } \bigwedge \Phi = \bigcup_{\varphi \in \Phi} \text{bnd } \varphi$, and $\text{bnd } \forall x \varphi = \text{bnd } \varphi \cup \{x\}$. Intuitively, $\text{bnd } \varphi$ is the set of variables x such that the prefix $\forall x$ occurs in φ . If $\text{bnd } \varphi = \emptyset$, then φ is *quantifier-free*.

Given $\varphi \in \mathcal{L}$, we define the set of *free variables* in φ , denoted by $\text{free } \varphi$, by formula recursion. Namely, $\text{free } \pi = \text{var } \pi$ if π is prime, $\text{free } \neg \varphi = \text{free } \varphi$, $\text{free } \bigwedge \Phi = \bigcup_{\varphi \in \Phi} \text{free } \varphi$, and $\text{free } \forall x \varphi = \text{free } \varphi \setminus \{x\}$. Intuitively, $\text{free } \varphi$ is the set of variables in φ that are not associated with a quantifier. For $X \subseteq \mathcal{L}$, we define $\text{free } X = \bigcup_{\varphi \in X} \text{free } \varphi$.

Strictly speaking, $\text{bnd } \varphi$ and $\text{free } \varphi$ need not be disjoint. For example, suppose

$$\varphi = (x \leq y \wedge \forall x \exists y x + y = 0).$$

Here, \leq is a binary relation symbol and 0 is a constant symbol. In this formula, the first occurrences of x and y are both free, whereas the others are bound. Hence, $\text{bnd } \varphi = \text{free } \varphi = \{x, y\}$. Once we move beyond the syntax of formulas and establish their logical relationships, we will see that φ is logically equivalent to

$$\varphi' = (x \leq y \wedge \forall u \exists v u + v = 0).$$

In this formula, $\text{bnd } \varphi' = \{u, v\}$ and $\text{free } \varphi' = \{x, y\}$, so that $\text{bnd } \varphi'$ and $\text{free } \varphi'$ are disjoint. In general, given any formula φ , there is a logically equivalent φ' with $\text{bnd } \varphi' \cap \text{free } \varphi' = \emptyset$. Hence, for most purposes, we may assume that no variable is both bound and free.

A *sentence*, or *closed formula*, is a formula φ such that $\text{free } \varphi = \emptyset$. The set of sentences is denoted by \mathcal{L}^0 . The set of finitary sentences is $\mathcal{L}_{\text{fin}}^0 = \mathcal{L}^0 \cap \mathcal{L}_{\text{fin}}$. Note that $\mathcal{L}_{\text{fin}}^0$ is the set of sentences used in first-order logic. An *open formula* is a formula that has one or more free variables.

If $x_1, \dots, x_n \in \text{Var}$ are distinct, we will write $\varphi = \varphi(x_1, \dots, x_n)$ or $\varphi = \varphi(\vec{x})$ to mean that $\text{free } \varphi \subseteq \{x_1, \dots, x_n\}$. Similarly, for $t \in \mathcal{T}$, we write $t = t(\vec{x})$ to mean that $\text{var } t \subseteq \{x_1, \dots, x_n\}$.

Definition 5.1.7. For $\varphi \in \mathcal{L}$, the set $\text{sym } \varphi \subseteq L$ is defined recursively as follows:

- (i) $\text{sym } s = t = \text{sym } s \cup \text{sym } t$,
- (ii) $\text{sym } r t_1 \cdots t_n = \{r\} \cup \text{sym } t_1 \cup \cdots \cup \text{sym } t_n$,
- (iii) $\text{sym } \neg \varphi = \text{sym } \varphi$,
- (iv) $\text{sym } \bigwedge \Phi = \bigcup_{\varphi \in \Phi} \text{sym } \varphi$, and
- (v) $\text{sym } \forall x \varphi = \text{sym } \varphi$.

Intuitively, $\text{sym } \varphi$ is the set of extralogical symbols occurring in φ . Note that $\text{sym } \varphi$, like $\text{var } \varphi$, is also countable. In addition, we define

$$\text{con } \varphi = \{c \in \text{sym } \varphi \mid c \text{ is a constant symbol}\},$$

which denotes the (countable) set of constant symbols occurring in φ .

5.1.5 Substitutions

(S:subs) A *substitution* is a function $\sigma : \text{Var} \rightarrow \mathcal{T}$. Such a function can be extended to $\sigma : \mathcal{T} \rightarrow \mathcal{T}$ by $c^\sigma = c$ and $(ft)^\sigma = ft_1^\sigma \cdots ft_n^\sigma$. Given $\varphi \in \mathcal{L}$, we want to define $\varphi^\sigma \in \mathcal{L}$ so that φ^σ denotes the result of substituting for every term t in φ the new term t^σ . We do this by transfinite recursion on the rank of φ .

Suppose $\varphi \in \mathcal{L}$ is prime. Then $\varphi = (s = t)$ or $\varphi = rt_1 \cdots t_n$. In the first case, we define $\varphi^\sigma = (s^\sigma = t^\sigma)$. In the second case, we define $\varphi^\sigma = rt_1^\sigma \cdots t_n^\sigma$. In this way, we have defined the map $\sigma \mapsto \varphi^\sigma$ for every $\varphi \in \mathcal{L}$ with $\text{rk } \varphi = 0$.

Let α be a nonzero ordinal and assume $\sigma \mapsto \varphi^\sigma$ has been defined for every $\varphi \in \mathcal{L}$ with $\text{rk } \varphi < \alpha$. Fix $\varphi \in \mathcal{L}$ with $\text{rk } \varphi = \alpha$. By the unique formula reconstruction property, one of the following holds:

- (i) $\varphi = \neg\psi$, where $\text{rk } \psi < \alpha$,
- (ii) $\varphi = \bigwedge \Phi$, where $\text{rk } \theta < \alpha$ for all $\theta \in \Phi$, or
- (iii) $\varphi = \forall x\psi$, where $\text{rk } \psi < \alpha$.

In the first two cases, define $\varphi^\sigma = \neg\psi^\sigma$ and $\varphi^\sigma = \bigwedge_{\theta \in \Phi} \theta^\sigma$, respectively. In the third case, define $\varphi^\sigma = \forall x\psi^\tau$, where $\tau : \text{Var} \rightarrow \mathcal{T}$ is the substitution defined by $x^\tau = x$ and $y^\tau = y^\sigma$ whenever $y \neq x$. By transfinite recursion on $\text{rk } \varphi$, this defines $\sigma \mapsto \varphi^\sigma$ for all $\varphi \in \mathcal{L}$.

If $x \in \text{Var}$ and $t \in \mathcal{T}$, we use the notation t/x to denote the substitution $\sigma : \text{Var} \rightarrow \mathcal{T}$ defined by $x^\sigma = t$ and $y^\sigma = y$ for $y \neq x$. We read t/x as “ t for x ” and write φ^σ as $\varphi(t/x)$. This is extended in the natural way for $\vec{x} = \langle x_1, \dots, x_n \rangle$ and $\vec{t} = \langle t_1, \dots, t_n \rangle$.

Note that, in general, $\varphi(t_1t_2/x_1x_2) \neq \varphi(t_1/x_1)(t_2/x_2)$. For example, if $\varphi = x_1 < x_2$, then $\varphi(x_2x_1/x_1x_2) = x_2 < x_1$, but $\varphi(x_2/x_1)(x_1/x_2) = x_1 < x_1$.

(P:var-sub-term) **Proposition 5.1.8.** *Let $s, t \in \mathcal{T}$ and $x \in \text{Var}$. Then*

$$\text{var } s(t/x) \subseteq (\text{var } s \setminus \{x\}) \cup \text{var } t.$$

Proof. We prove this by induction on s . If $s = x$, then $s(t/x) = t$, so $\text{var } s(t/x) = \text{var } t$, and the result holds. If $s = y$, where $y \neq x$, or $s = c$, then $s(t/x) = s$ and $x \notin \text{var } s$. Hence, $\text{var } s(t/x) = \text{var } s = \text{var } s \setminus \{x\}$, and the result holds.

Now suppose the result holds for t_1, \dots, t_n , and let $s = ft_1 \cdots t_n$. Let $t'_i = t_i(t/x)$, so that $\text{var } t'_i \subseteq (\text{var } t_i \setminus \{x\}) \cup \text{var } t$ and $s(t/x) = ft'_1 \cdots t'_n$. Then

$$\begin{aligned} \text{var } s(t/x) &= \bigcup_{i=1}^n \text{var } t'_i \\ &\subseteq \bigcup_{i=1}^n (\text{var } t_i \setminus \{x\}) \cup \text{var } t \\ &= ((\bigcup_{i=1}^n \text{var } t_i) \setminus \{x\}) \cup \text{var } t \\ &= (\text{var } s \setminus \{x\}) \cup \text{var } t, \end{aligned}$$

and the result holds. □

(P:free-sub-form) **Proposition 5.1.9.** *Let $\varphi \in \mathcal{L}$, $x \in \text{Var}$, and $t \in \mathcal{T}$. Then*

$$\text{free } \varphi(t/x) \subseteq (\text{free } \varphi \setminus \{x\}) \cup \text{var } t.$$

Proof. We prove this by induction on φ . An argument like the one in the proof of Proposition 5.1.8 covers the cases where φ is prime, $\varphi = \neg\psi$, and $\varphi = \bigwedge \Phi$. Suppose $\varphi = \forall x\psi$. Then $\text{free } \varphi = \text{free } \psi \setminus \{x\}$. In particular, $x \notin \text{free } \varphi$. Moreover, $\varphi(t/x) = \varphi$, so that $\text{free } \varphi(t/x) = \text{free } \varphi = \text{free } \varphi \setminus \{x\}$, and the result holds.

Now suppose $\varphi = \forall y\psi$, where $y \neq x$. Then $\varphi(t/x) = \forall y\psi(t/x)$, and $\text{free } \varphi(t/x) = \text{free } \psi(t/x) \setminus \{y\}$. Assume $z \in \text{free } \varphi(t/x)$, but $z \notin (\text{free } \varphi \setminus \{x\}) \cup \text{var } t$. Then $z \in \text{free } \psi(t/x)$ and $z \neq y$. By the inductive hypothesis, the result holds for ψ . Hence, $z \in (\text{free } \psi \setminus \{x\}) \cup \text{var } t$. But $z \notin \text{var } t$. Thus, $z \in \text{free } \psi$ and $z \neq x$. Since $\text{free } \varphi = \text{free } \psi \setminus \{y\}$, we have $z \in \text{free } \varphi$. It follows that $z \in \text{free } \varphi \setminus \{x\}$, a contradiction. \square

Corollary 5.1.10. *If $\varphi(x) \in \mathcal{L}$ and t is a ground term, then $\varphi(t) \in \mathcal{L}^0$.*

Proof. Let $\varphi(x) \in \mathcal{L}$ and let t be a ground term. Then $\text{free } \varphi \subseteq \{x\}$ and $\text{var } t = \emptyset$. By Proposition 5.1.9, we have $\text{free } \varphi(t) = \emptyset$, so that $\varphi(t)$ is a sentence. \square

Using substitutions, we introduce the shorthand,

$$\exists!x\varphi = \exists x\varphi \wedge \forall xy(\varphi \wedge \varphi(y/x) \rightarrow x = y),$$

where $y \notin \text{var } \varphi$.

5.2 Predicate calculus

(S:pred-calc) In this section, we define both the deductive and inductive derivability relations. As in the propositional case, we will denote them both by \vdash . We begin with the deductive case. As described in Section 5.2.6, the inductive case will require no modification from its presentation in Chapter 3.

For deductive derivability, we wish to define a relation \vdash from $\mathfrak{P}\mathcal{L}^0$ to \mathcal{L}^0 such that $X \vdash \varphi$ captures what it means to say a sentence φ can be logically deduced from the sentences in X . Our aim here is to do this through natural deduction, as we did in Section 3.1.2 for the propositional language \mathcal{F} . We will keep all the rules in Definition 3.1.3, and add two rules each for \forall and $=$. Ideally, we would like our new rules to be the following:

1. if $X \vdash \forall x\varphi(x)$ and t is a ground term, then $X \vdash \varphi(t)$,
2. if $c \notin \text{con}(X \cup \varphi(x))$ and $X \vdash \varphi(c)$, then $X \vdash \forall x\varphi(x)$,
3. $\vdash t = t$ for all ground terms t , and
4. if $X \vdash s = t, \varphi(s)$, then $X \vdash \varphi(t)$.

The problem with these rules is (2). We may not have enough constants in our language to ensure there exists a $c \notin \text{con}(X \cup \varphi(x))$. As shown later in Proposition 5.2.16, we can always add constants to our language without affecting derivability. But until that can be established, these rules will not be easy to work with.

We therefore take a slightly different approach. We define the derivability relation from $\mathfrak{P}\mathcal{L}$ to \mathcal{L} . That is, we allow ourselves to use open formulas in our derivations. In an open formula, we will treat free variables like constants. In this way, every language will effectively have an uncountable number of constants available for use in (2). This still doesn't fully resolve the problem, since X itself can be uncountable. But in Theorem 5.2.11, we will prove σ -compactness, so that we need only consider countable X .

Reasoning with sentences, however, is our primary concern in the bulk of what we want to do. Both deductive and inductive theories consist entirely of sentences. As such, after presenting our system of natural deduction and proving σ -compactness, we will look at expanding our language by adding additional constants. This will give us two ways to connect our natural deduction to reasoning with sentences. These are presented in Propositions 5.2.17 and 5.2.18.

We then define deductive and inductive theories, and finally finish the section with a presentation of a Hilbert-type calculus. This is the calculus used by Karp in [16], and we will need it in order to apply her completeness result in the setting of predicate logic.

By allowing open formulas, however, we introduce a new problem. Suppose $\mathcal{L} = \mathcal{L}\{\circ, e\}$, where \circ is a binary operation symbol and e is a constant symbol. Let $\varphi(x) = \exists y x \neq y$ and $t = y \circ e$. Then $\varphi(t) = \exists y y \circ e \neq y$. If we interpret these formulas in group theory, where e is the group identity and \circ is the group operation, then $\forall x \varphi(x) = \forall x \exists y x \neq y$ is a true sentence in every group that has more than one element. And yet, the sentence $\varphi(t)$ is always false in that context. Hence, $\forall x \varphi(x)$ cannot logically imply $\varphi(t)$. This would not violate (1), because t is not a ground term. But when we remove that restriction on t , this will become a problem. The issue is that the variable $y \in \text{var } t$ is a bound variable in φ , so after the substitution, it becomes bound. If free variables are to be treated as constants, then variables inside terms must become free after a substitution. To ensure this, we will need to avoid substitutions that “collide” with bound variables.

5.2.1 Free substitutions

Let $\varphi \in \mathcal{L}$, $\zeta \in \text{Sf } \varphi$, and $x \in \text{Var}$. We say that ζ is in the scope of $\forall x$ in φ if there exists ψ such that $\forall x \psi \in \text{Sf } \varphi$ and $\zeta \in \text{Sf } \psi$. We say that ζ has a free occurrence of x in φ if $x \in \text{free } \zeta$ and ζ is not in the scope of $\forall x$ in φ . Note that if $\varphi' \in \text{Sf } \varphi$, $\zeta \in \text{Sf } \varphi'$, and ζ has a free occurrence of x in φ , then ζ has a free occurrence of x in φ' .

(P:free-occurrence)

Proposition 5.2.1. *If ζ has a free occurrence of x in φ , then $x \in \text{free } \varphi$.*

Proof. If $\zeta = \varphi$, the result is immediate, so assume $\zeta \neq \varphi$. Then φ is not prime, so we may write $\varphi = \neg\varphi'$, $\varphi = \bigwedge \Phi$, or $\varphi = \forall y\varphi'$. By induction on φ , we may assume the result is true for φ' and for all $\theta \in \Phi$.

In the first case, $\zeta \in \text{Sf } \varphi'$, so ζ has a free occurrence of x in φ' . Hence, $x \in \text{free } \varphi' = \text{free } \varphi$. In the second case, $\zeta \in \text{Sf } \theta$ for some $\theta \in \Phi$. Thus, $x \in \text{free } \theta \subseteq \text{free } \varphi$. Similarly, in the third case, we get $x \in \text{free } \varphi'$. But ζ is not in the scope of $\forall x$ in φ , so $y \neq x$. Hence, $\text{free } \varphi = \text{free } \varphi'$. \square

Let $x, y \in \text{Var}$ and $\varphi \in \mathcal{L}$. We say that y is not free for x in φ if there exists $\zeta \in \text{Sf } \varphi$ such ζ is in the scope of $\forall y$ in φ and ζ has a free occurrence of x in φ . Otherwise, y is free for x in φ . For $t \in \mathcal{T}$, we say that t is free for x in φ if y is free for x in φ for all $y \in \text{var } t$. More generally, a substitution σ is free for x in φ if x^σ is free for x in φ , for all $x \in \text{Var}$.

Proposition 5.2.2. *If $\text{bnd } \varphi \cap (\text{var } t \setminus \{x\}) = \emptyset$, then t is free for x in φ . In particular, y is free for x in φ if $y = x$ or $y \notin \text{bnd } \varphi$.*

Proof. Suppose t is not free for x in φ . Then there exists $y \in \text{var } t$ such that a free occurrence of x occurs inside the scope of $\forall y$. In particular, we must have $y \neq x$ and $y \in \text{bnd } \varphi$, so that $y \in \text{bnd } \varphi \cap (\text{var } t \setminus \{x\})$. \square

(P:sub-free-only) **Proposition 5.2.3.** *Let $y \notin \text{var } \varphi$ and $\zeta \in \text{Sf } \varphi(y/x)$. If ζ has a free occurrence of y in $\varphi(y/x)$, then ζ is not in the scope of $\forall x$ in $\varphi(y/x)$.*

Proof. The proof is by induction on φ and follows the same lines as the proof of Proposition 5.2.1. \square

(C:free-sub-only) **Corollary 5.2.4.** *If $y \notin \text{var } \varphi$, then x is free for y in $\varphi(y/x)$.*

Proof. Suppose x is not free for y in $\varphi(y/x)$. Then there exists $\zeta \in \text{Sf } \varphi(y/x)$ such that ζ is in the scope of $\forall x$ in $\varphi(y/x)$ and ζ has a free occurrence of y in $\varphi(y/x)$. But this contradicts Proposition 5.2.3. \square

5.2.2 Natural deduction

(D:pred-derivability) **Definition 5.2.5.** The *derivability relation*, denoted by \vdash or $\vdash_{\mathcal{L}}$, is the smallest relation from $\mathfrak{P}\mathcal{L}$ to \mathcal{L} satisfying (i)–(vi) in Definition 3.1.3, as well as the following:

- (vii) if $X \vdash \forall x\varphi$, then $X \vdash \varphi(t/x)$ when t is free for x in φ ,
- (viii) if $x \notin \text{free } X$ and $X \vdash \varphi$, then $X \vdash \forall x\varphi$,
- (ix) $\vdash t = t$ for all $t \in \mathcal{T}$, and
- (x) if $X \vdash s = t, \varphi(s/x)$, then $X \vdash \varphi(t/x)$ when s and t are free for x in φ .

Since x is always free for x in φ , (vii) implies

- (vii)' if $X \vdash \forall x\varphi$, then $X \vdash \varphi$.

^(R:mon) **Remark 5.2.6.** If $X \vdash_{\mathcal{L}} \varphi$ and $\mathcal{L} \subseteq \mathcal{L}'$, then $X \vdash_{\mathcal{L}'} \varphi$. To see this, let $\mathcal{L} \subseteq \mathcal{L}'$ and define $\vdash' = \vdash_{\mathcal{L}'} \cap (\mathfrak{P} \mathcal{L} \times \mathcal{L})$. Since \vdash' satisfies (i)–(x) for \mathcal{L} , we have $\vdash_{\mathcal{L}} \subseteq \vdash'$.

^(R:pred-fin-vs-infin) **Remark 5.2.7.** The *finitary derivability relation* is the smallest relation \vdash_{fin} from $\mathfrak{P} \mathcal{L}_{\text{fin}}$ to \mathcal{L}_{fin} such that conditions (i)–(x) from Definition 5.2.5 hold, with the exception that in (iii) and (iv), we require Φ to be finite. The finitary derivability relation is a typical natural-deduction calculus for first-order logic. Clearly, $\vdash_{\text{fin}} \subseteq \vdash$. As we will see in Proposition 5.3.15, if $X \subseteq \mathcal{L}_{\text{fin}}$, $\varphi \in \mathcal{L}_{\text{fin}}$, and $X \vdash \varphi$, then $X \vdash_{\text{fin}} \varphi$. In other words, when restricted to finitary formulas, infinitary calculus cannot produce any new inferences beyond those already available in first-order logic.

The proof of Proposition 3.1.8, which is based on (i)–(vi), is still valid here. Throughout the rest of this chapter, unless otherwise indicated, we will use lowercase Roman numerals refer to Definition 5.2.5 and letters refer to Proposition 3.1.8.

^(P:bind-rename) **Proposition 5.2.8 (Bound renaming).** For any $\varphi \in \mathcal{L}$ and $y \notin \text{var } \varphi$, we have $\forall x \varphi \vdash \forall y \varphi(y/x)$ and $\forall y \varphi(y/x) \vdash \forall x \varphi$.

Proof. If $y = x$, the result follows from (i). Assume, then, that $y \neq x$. Since $y \notin \text{var } \varphi$, it follows from (vii) that $\forall x \varphi \vdash \varphi(y/x)$. Hence, by (viii), we have $\forall x \varphi \vdash \forall y \varphi(y/x)$.

Let $\varphi' = \varphi(y/x)$. Corollary 5.2.4 implies that x is free for y in φ' . Hence, by (vii), we have $\forall y \varphi' \vdash \varphi'(x/y) = \varphi$. But $x \notin \text{free } \forall y \varphi'$, so (viii) implies $\forall y \varphi' \vdash \forall x \varphi$. \square

Remark 5.2.9. Bound renaming gives us the following alternate to (viii):

(viii)' if $y \notin \text{free } X \cup \text{var } \varphi$ and $X \vdash \varphi(y/x)$, then $X \vdash \forall x \varphi$,

To see this, suppose $y \notin \text{free } X \cup \text{var } \varphi$ and $X \vdash \varphi(y/x)$. Then (viii) implies $X \vdash \forall y \varphi(y/x)$ and Proposition 5.2.8 gives $\forall y \varphi(y/x) \vdash \forall x \varphi$.

We now prove σ -compactness. In the predicate case, the theorem is stronger, in that we can not only pass to a countable subset of formulas. We can also pass to a countable subset of extralogical symbols. We begin with the basic version, which is the analogue of the propositional version.

Proposition 5.2.10. Let $X \subseteq \mathcal{L}$ and $\varphi \in \mathcal{L}$. Then $X \vdash \varphi$ if and only if there exists a countable subset $X_0 \subseteq X$ such that $X_0 \vdash \varphi$.

Proof. As in the proof of Theorem 3.1.10, we will prove that the following are equivalent:

$$X \vdash \varphi, \tag{5.2.1} \boxed{\text{pred-sig-cpctness-1}}$$

$$\text{there exists countable } X_0 \subseteq X \text{ such that } \bigwedge X_0 \vdash \varphi, \text{ and} \tag{5.2.2} \boxed{\text{pred-sig-cpctness-2}}$$

$$\text{there exists countable } X_0 \subseteq X \text{ such that } X_0 \vdash \varphi. \tag{5.2.3} \boxed{\text{pred-sig-cpctness-3}}$$

The proof of Theorem 3.1.10 carries through to show that (5.2.2) implies (5.2.3), and (5.2.3) implies (5.2.1). Define \vdash' so that $X \vdash' \varphi$ if and only if $X \vdash \varphi$ and (5.2.2) holds. Since Lemma 3.1.9 is still valid for \vdash , the proof of Theorem 3.1.10 shows that (i)–(iv) hold for \vdash' . It is straightforward to verify that \vdash' satisfies (vii)–(x). \square

(T:pred-sig-cpctness) **Theorem 5.2.11 (σ -compactness).** *Let $X \subseteq \mathcal{L}$ and $\varphi \in \mathcal{L}$. Then $X \vdash_{\mathcal{L}} \varphi$ if and only if there exist countable $X_0 \subseteq X$ and $L_0 \subseteq L$ such that $X_0 \vdash_{\mathcal{L}_0} \varphi$.*

Proof. The if direction follows from (ii) and Remark 5.2.6. For the only if direction, define $\vdash' \subseteq \mathfrak{P}\mathcal{L} \times \mathcal{L}$ by $X \vdash' \varphi$ if $X_0 \vdash_{\mathcal{L}_0} \varphi$ for some countable $X_0 \subseteq X$ and $L_0 \subseteq L$. It suffices to show that \vdash' satisfies (i)–(x).

Suppose $\varphi \in \mathcal{L}$. Let $L_0 = \text{sym } \varphi$. Then $L_0 \subseteq L$ is countable, and by (i) for \mathcal{L}_0 , we have $\varphi \vdash_{\mathcal{L}_0} \varphi$. Thus, $\varphi \vdash' \varphi$, and \vdash' satisfies (i).

Let $\Phi \subseteq \mathcal{L}$ be countable and suppose $X \vdash' \theta$ for all $\theta \in \Phi$. For each $\theta \in \Phi$, choose countable $X_\theta \subseteq X$ and $L_\theta \subseteq L$ such that $X_\theta \vdash_{\mathcal{L}_\theta} \theta$. Let $X_0 = \bigcup_{\theta \in \Phi} X_\theta$ and $L_0 = \bigcup_{\theta \in \Phi} L_\theta$, both of which are countable. Remark 5.2.6 implies $X_\theta \vdash_{\mathcal{L}_0} \theta$ for all $\theta \in \Phi$. Hence, by (ii) for \mathcal{L}_0 , we have $X_0 \vdash_{\mathcal{L}_0} \theta$ for all $\theta \in \Phi$. Therefore, by (iv) for \mathcal{L}_0 , it follows that $X_0 \vdash_{\mathcal{L}_0} \bigwedge \Phi$. Thus, $X \vdash' \bigwedge \Phi$, and \vdash' satisfies (iv).

The proofs of (ii), (iii), and (v)–(x) are similar. \square

Since they are based on (i)–(vi), Propositions 3.1.11 and 3.1.12 hold here as well.

We finish this subsection with a result that we will need later. A *variable permutation* is a bijection $\pi : \text{Var} \rightarrow \text{Var}$. We extend a variable permutation to $\pi : \mathcal{T} \rightarrow \mathcal{T}$ by $c^\pi = c$ and $(ft_1 \cdots t_n)^\pi = ft_1^\pi \cdots t_n^\pi$. For $\varphi \in \mathcal{L}$, we define φ^π by $(s = t)^\pi = (s^\pi = t^\pi)$, $(rt_1 \cdots t_n)^\pi = rt_1^\pi \cdots t_n^\pi$, $(\neg\varphi)^\pi = \neg\varphi^\pi$, $(\bigwedge \Phi)^\pi = \bigwedge_{\varphi \in \Phi} \varphi^\pi$, and $(\forall x\varphi)^\pi = \forall x^\pi \varphi^\pi$.

(P:perm-vars) **Proposition 5.2.12.** *If π is a variable permutation, then $X \vdash \varphi$ if and only if $X^\pi \vdash \varphi^\pi$.*

Proof. Define \vdash' by $X \vdash' \varphi$ if $X \vdash \varphi$ and $X^\pi \vdash \varphi^\pi$. It is straightforward to verify that \vdash' satisfies (i)–(x). Hence, $X \vdash \varphi$ implies $X^\pi \vdash \varphi^\pi$. Applying this result to π^{-1} gives the converse. \square

5.2.3 Constant expansions

We now connect our deductive system back to sentences. To do this, we will need to expand our language by adding additional constant symbols.

Let L be an extralogical signature with corresponding language \mathcal{L} . Let C be a set of constant symbols. Then $\mathcal{L}C$ denotes the language corresponding to the extralogical signature $L \cup C$. If $C = \{c\}$, then we write $\mathcal{L}c$ for $\mathcal{L}C$. The language $\mathcal{L}C$ is called a *constant expansion* of \mathcal{L} .

Let C be a countable set of constant symbols. A *C-substitution* is an injective function $\underline{\sigma} : C \rightarrow \text{Var}$. Given a C -substitution, we extend it to $\underline{\sigma} : \mathcal{T}_{\mathcal{L}C} \rightarrow \text{Var}$ by $x^\underline{\sigma} = x$, $c^\underline{\sigma} = \underline{\sigma}(c)$ for $c \in C$, and $(ft_1 \cdots t_n)^\underline{\sigma} = ft_1^\underline{\sigma} \cdots t_n^\underline{\sigma}$. Define $\varphi^\underline{\sigma}$

recursively by $(s = t)^\sigma = (s^\sigma = t^\sigma)$, $(rt_1 \cdots t_n)^\sigma = rt_1^\sigma \cdots t_n^\sigma$, $(\neg\varphi)^\sigma = \neg\varphi^\sigma$, $(\bigwedge \Phi)^\sigma = \bigwedge_{\varphi \in \Phi} \varphi^\sigma$, and $(\forall x\varphi)^\sigma = \forall x\varphi^\sigma$. Intuitively, t^σ and φ^σ are obtained from t and φ , respectively, by replacing each occurrence of $c \in C$ with c^σ . For $X \subseteq \mathcal{L}$, we let $X^\sigma = \{\varphi^\sigma \mid \varphi \in X\}$. Note that if $X \subseteq \mathcal{LC}$, then $X^\sigma \subseteq \mathcal{L}$. For $V \subseteq \mathbf{Var}$, we say that $\underline{\sigma}$ *avoids* V if $c^\sigma \notin V$ for all $c \in C$. Since \mathbf{Var} is uncountable, given any countable $V_0 \subseteq \mathbf{Var}$, there exists a C -substitution that avoids V_0 . By term induction, it is easy to see that $s(t/x)^\sigma = s^\sigma(t^\sigma/x)$ when $\underline{\sigma}$ avoids x . Using this and formula induction, we have $\varphi(t/x)^\sigma = \varphi^\sigma(t^\sigma/x)$.

If c is a constant symbol and $y \in \mathbf{Var}$, we use the notation y/c to denote the $\{c\}$ -substitution $\underline{\sigma}$ defined by $c^\sigma = y$. We read y/c as “ y for c ” and write φ^σ as $\varphi(y/c)$. In this case, the identity $\varphi(t/x)^\sigma = \varphi^\sigma(t^\sigma/x)$ becomes $\varphi(t/x)(y/c) = \varphi'(t'/x)$, where $t' = t(y/c)$ and $\varphi' = \varphi(y/c)$. This is extended in the natural way for $\vec{c} = \langle c_1, \dots, c_n \rangle$ and $\vec{x} = \langle x_1, \dots, x_n \rangle$, provided the c_i 's and the x_i 's are distinct.

(P:C-sub-taut) **Proposition 5.2.13.** *If $\underline{\sigma}$ and $\underline{\sigma}'$ both avoid $\text{var } \varphi$, then there is a variable permutation such that $\varphi^{\underline{\sigma}'} = (\varphi^\sigma)^\pi$. In particular, $\vdash_{\mathcal{L}} \varphi^\sigma$ if and only if $\vdash_{\mathcal{L}} \varphi^{\underline{\sigma}'}$.*

Proof. Since C is countable and both $\underline{\sigma}$ and $\underline{\sigma}'$ are injective, we may choose a bijection $\pi : \mathbf{Var} \rightarrow \mathbf{Var}$ such that $\pi \circ \underline{\sigma} = \underline{\sigma}'$ and $x^\pi = x$ for all x outside the ranges of $\underline{\sigma}$ and $\underline{\sigma}'$. We prove that $\varphi^{\underline{\sigma}'} = (\varphi^\sigma)^\pi$ by induction on φ . The proof is entirely straightforward except for the inductive step $\varphi = \forall x\psi$. For this, let $\varphi = \forall x\psi$ and assume $\underline{\sigma}$ and $\underline{\sigma}'$ both avoid $\text{var } \varphi$. Then $\underline{\sigma}$ and $\underline{\sigma}'$ both avoid $\text{var } \psi$, so by the inductive hypothesis, $\psi^{\underline{\sigma}'} = (\psi^\sigma)^\pi$. Since $\underline{\sigma}$ and $\underline{\sigma}'$ avoid $\text{var } \varphi$, we have $x^\pi = x$. Thus,

$$(\varphi^\sigma)^\pi = (\forall x\psi^\sigma)^\pi = \forall x^\pi(\psi^\sigma)^\pi = \forall x\psi^{\underline{\sigma}'} = \varphi^{\underline{\sigma}'}$$

For the final result, we apply Proposition 5.2.12. □

(P:C-sub-deriv) **Proposition 5.2.14.** *Let $X \subseteq \mathcal{LC}$ and $\varphi \in \mathcal{LC}$. If $\underline{\sigma}$ and $\underline{\sigma}'$ are C -substitutions that both avoid $\text{var } X \cup \{\varphi\}$, then $X^\sigma \vdash_{\mathcal{L}} \varphi^\sigma$ if and only if $X^{\underline{\sigma}'} \vdash_{\mathcal{L}} \varphi^{\underline{\sigma}'}$.*

Proof. Suppose $X^\sigma \vdash_{\mathcal{L}} \varphi^\sigma$. Choose countable $X_0 \subseteq X$ such that $X_0^\sigma \vdash_{\mathcal{L}} \varphi^\sigma$. Let $\psi = \bigwedge X_0$. Since $\psi^\sigma = \bigwedge_{\theta \in X_0} \theta^\sigma$ and $(\psi \rightarrow \varphi)^\sigma = \psi^\sigma \rightarrow \varphi^\sigma$, we have $\vdash_{\mathcal{L}} (\psi \rightarrow \varphi)^\sigma$. But $\underline{\sigma}$ and $\underline{\sigma}'$ both avoid $\text{var}(\psi \rightarrow \varphi)$. Hence, by Proposition 5.2.13, we have $\vdash_{\mathcal{L}} (\psi \rightarrow \varphi)^{\underline{\sigma}'}$, which implies $X^{\underline{\sigma}'} \vdash_{\mathcal{L}} \varphi^{\underline{\sigma}'}$. Reversing the roles of $\underline{\sigma}$ and $\underline{\sigma}'$ gives the converse. □

(P:con-elim) **Proposition 5.2.15 (Constant elimination).** *Let C be a countable set of constant symbols and suppose $X \vdash_{\mathcal{LC}} \varphi$. Then there exists a countable set $X_0 \subseteq X$ such that $X_0 \vdash_{\mathcal{LC}} \varphi$ and $X_0^\sigma \vdash_{\mathcal{L}} \varphi^\sigma$ whenever $\underline{\sigma}$ is a C -substitution that avoids $\text{var } X_0 \cup \{\varphi\}$. In particular, if $\underline{\sigma}$ avoids $\text{var } X \cup \{\varphi\}$, then $X^\sigma \vdash_{\mathcal{L}} \varphi^\sigma$.*

Proof. Define $\vdash' \subseteq \mathfrak{P}\mathcal{LC} \times \mathcal{LC}$ by $X \vdash' \varphi$ if $X \vdash_{\mathcal{LC}} \varphi$ and there exists countable $X_0 \subseteq X$ such that $X_0 \vdash_{\mathcal{LC}} \varphi$ and $X_0^\sigma \vdash_{\mathcal{L}} \varphi^\sigma$ whenever $\underline{\sigma}$ avoids $\text{var } X_0 \cup \{\varphi\}$. It suffices to show that \vdash' satisfies (i)–(x).

It is immediate that (i) and (ii) hold. For (iii), suppose $X \vdash' \bigwedge \Phi$. Choose countable $X_0 \subseteq X$ such that $X_0 \vdash_{\mathcal{LC}} \bigwedge \Phi$ and $X_0^\sigma \vdash_{\mathcal{L}} (\bigwedge \Phi)^\sigma = \bigwedge_{\theta \in \Phi} \theta^\sigma$

whenever $\underline{\sigma}$ avoids $\text{var } X_0 \cup \{\bigwedge \Phi\} = \text{var } X_0 \cup \Phi$. Now fix $\theta \in \Phi$. Then $X_0 \vdash_{\mathcal{L}C} \theta$ and $X_0^\sigma \vdash_{\mathcal{L}} \theta^\sigma$ whenever $\underline{\sigma}$ avoids $\text{var } X_0 \cup \Phi$. Suppose $\underline{\sigma}$ avoids $\text{var } X_0 \cup \{\theta\}$. Since $\text{var } X_0 \cup \Phi$ is countable, we may choose $\underline{\sigma}'$ that avoids $\text{var } X_0 \cup \Phi$. Since $\underline{\sigma}$ and $\underline{\sigma}'$ both avoid $\text{var } X_0 \cup \{\theta\}$, the inductive hypothesis and Proposition 5.2.14 give $X_0^\sigma \vdash_{\mathcal{L}} \theta^\sigma$. Therefore, $X \vdash' \theta$ and \vdash' satisfies (iii). The proofs for (iv)–(vi) are similar.

For (vii), suppose $X \vdash' \forall x\varphi$ and let t be free for x in φ . Choose countable $X_0 \subseteq X$ such that $X_0 \vdash \forall x\varphi$ and $X_0^\sigma \vdash (\forall x\varphi)^\sigma = \forall x\varphi^\sigma$ whenever $\underline{\sigma}$ avoids $\text{var } X \cup \{\forall x\varphi\}$. Now let $\underline{\sigma}$ avoid $\text{var } X \cup \{\varphi(t/x)\}$. As above, by Proposition 5.2.14, we may assume $\underline{\sigma}$ also avoids both $\text{var } X_0 \cup \{\forall x\varphi\}$ and $\text{var } t$. Then by hypothesis, $X_0^\sigma \vdash (\forall x\varphi)^\sigma = \forall x\varphi^\sigma$.

We now show that t^σ is free for x in φ^σ . Suppose not. Then there exists $y \in \text{var } t^\sigma$ such that a free occurrence of x in φ^σ occurs within the scope of $\forall y$ in φ^σ . But $\underline{\sigma}$ avoids x , so every occurrence of x in φ^σ is an occurrence of x in φ . Also, by the definition of C -substitutions, any occurrence of the quantifier $\forall y$ in φ^σ is an occurrence of $\forall y$ in φ . Hence, there is a free occurrence of x in φ that occurs within the scope of $\forall y$ in φ . Since t is free for x in φ , we must have $y \notin \text{var } t$. On the other hand, since $\forall y$ occurs in φ , we have $y \in \text{var } \varphi$. Therefore, $\underline{\sigma}$ avoids y . It follows that $y \notin \text{var } t^\sigma$, a contradiction.

Since t^σ is free for x in φ^σ and $X_0^\sigma \vdash \forall x\varphi^\sigma$, it follows from (vii) that $X_0^\sigma \vdash \varphi^\sigma(t^\sigma/x)$. But $\varphi^\sigma(t^\sigma/x) = \varphi(t/x)^\sigma$, so that $X_0^\sigma \vdash \varphi(t/x)^\sigma$, showing that \vdash' satisfies (vii). The proofs of (viii)–(x) are the similar. \square

(P:add-constants) **Proposition 5.2.16.** *Let $\mathcal{L}C$ be a constant expansion of \mathcal{L} . Let $X \subseteq \mathcal{L}$ and $\varphi \in \mathcal{L}$. Then $X \vdash_{\mathcal{L}} \varphi$ if and only if $X \vdash_{\mathcal{L}C} \varphi$.*

Proof. The only if direction follows from Remark 5.2.6. For the if direction, suppose $X \vdash_{\mathcal{L}C} \varphi$. By Theorem 5.2.11, we may choose countable $X_0 \subseteq X$ and $L_0 \subseteq L \cup C$ such that $X_0 \vdash_{\mathcal{L}C_0} \varphi$. Let $C_0 = (L_0 \cap C) \setminus L$ denote the constant symbols that are in L_0 but not in L . Then C_0 is countable and $\mathcal{L}_0 \subseteq \mathcal{L}C_0$. Hence, by Remark 5.2.6, we have $X_0 \vdash_{\mathcal{L}C_0} \varphi$. Let $\underline{\sigma}$ be a C_0 -substitution that avoids $\text{var } X_0 \cup \{\varphi\}$. Proposition 5.2.15 gives us $X_0^\sigma \vdash_{\mathcal{L}} \varphi^\sigma$. But $\text{sym } \varphi \subseteq L$, so $C_0 \cap \text{sym } \varphi = \emptyset$. In other words, none of the constants in C_0 appear in φ . By induction on φ , it follows that $\varphi^\sigma = \varphi$. Similarly, $X_0^\sigma = X_0$. Therefore, $X_0 \vdash_{\mathcal{L}} \varphi$, which gives $X \vdash_{\mathcal{L}} \varphi$. \square

5.2.4 Deduction with sentences

Using constant expansions, we can connect derivability back to sentences in two ways. The first is to verify the four rules given at the beginning of this section.

(P:pred-derivability) **Proposition 5.2.17.** *The derivability relation in \mathcal{L} , when restricted to sentences, satisfies (i)–(vi) in Definition 3.1.3, as well as the following:*

(vii)⁰ if $X \vdash \forall x\varphi(x)$ and t is a ground term, then $X \vdash \varphi(t)$,

(viii)⁰ if $c \notin \text{con}(X \cup \varphi(x))$ and $X \vdash_{\mathcal{L}c} \varphi(c)$, then $X \vdash \forall x\varphi(x)$, and

(ix)⁰ $\vdash t = t$ for all ground terms t ,

(x)⁰ if $X \vdash s = t, \varphi(s)$, then $X \vdash \varphi(t)$.

Proof. Since \vdash satisfies (i)–(vi) and \mathcal{L}^0 is closed under negation and conjunction, \vdash still satisfies (i)–(vi) when restricted to sentences. Suppose $X \vdash \forall x\varphi(x)$ and t is a ground term. Then $\varphi(t)$ is a sentence, and (vii) implies (vii)⁰.

For (viii)⁰, suppose $c \notin \text{con}(X \cup \varphi(x))$ and $X \vdash_{\mathcal{L}c} \varphi(c)$. Since $\forall x\varphi(x)$ is a sentence, we only need to show that $X \vdash \forall x\varphi$. By Theorem 5.2.11, we may choose countable $X_0 \subseteq X$ such that $X_0 \vdash_{\mathcal{L}c} \varphi(c)$. Choose $y \notin \text{var } X_0 \cup \{\varphi\}$. Then the $\{c\}$ -substitution y/c avoids $\text{var } X_0 \cup \{\varphi(c)\}$. Hence, by Proposition 5.2.15, we have $X_0(y/c) \vdash \varphi(c)(y/c)$. But $\varphi(c)(y/c) = \varphi'(c')$, where $c' = c(y/c) = y$ and $\varphi' = \varphi(y/c) = \varphi$, since $c \notin \text{con } \varphi$. Thus, $\varphi(c)(y/c) = \varphi$. Similarly, since $c \notin \text{con } X_0$, we also have $X_0(y/c) = X_0$. Therefore, $X_0 \vdash \varphi(y)$. Now $y \notin \text{free } X_0 \cup \text{var } \varphi$. Hence, from (viii)', it follows that $X_0 \vdash \forall x\varphi$, which gives $X \vdash \forall x\varphi$.

If t is a ground term, then $t = t$ is a sentence, so (ix) implies (ix)⁰. Now suppose $X \vdash s = t, \varphi(s)$. Then $s = t$ is a sentence, which implies s and t are ground terms. Thus, $\varphi(t)$ is a sentence, and (x) implies (x)⁰. \square

The second way to connect derivability to sentences is to replace the free variables in open formulas with constants. For each $x \in \text{Var}$, choose a distinct constant symbol c_x that is not already in L . Let $C = \{c_x \mid x \in \text{Var}\}$. A *free eliminator* is a substitution $\sigma : \text{Var} \rightarrow \mathcal{T}_{\mathcal{L}C}$ such that for all $x \in \text{Var}$, either $x^\sigma = x$ or $x^\sigma = c_x$. If $x^\sigma = c_x$ for all $x \in \text{Var}$, then σ is called a *full free eliminator*. Note that if σ is a full free eliminator, then φ^σ is a sentence for all $\varphi \in \mathcal{L}$.

(P:free-elim-deriv) **Proposition 5.2.18.** *Let $X \subseteq \mathcal{L}$ and $\varphi \in \mathcal{L}$, and let σ be a free eliminator. Then $X \vdash_{\mathcal{L}} \varphi$ if and only if $X^\sigma \vdash_{\mathcal{L}C} \varphi^\sigma$.*

Proof. Suppose $X \vdash_{\mathcal{L}} \varphi$. Let $V' = \{x \in \text{Var} \mid x^\sigma \neq x\}$ and let $C' = \{c_x \mid x \in V'\}$. Then $X \vdash_{\mathcal{L}C'} \varphi$ by Proposition 5.2.16. By Theorem 5.2.11, there exists countable $C_0 \subseteq C'$ such that $X \vdash_{\mathcal{L}C_0} \varphi$. Define the C_0 -substitution $\underline{\sigma}$ by $c_x^{\underline{\sigma}} = x$, so that $(\psi^{\underline{\sigma}})^\sigma = \psi$ for all $\psi \in \mathcal{L}C_0$, and $(\psi^\sigma)^\sigma = \psi$ whenever $\psi \in \mathcal{L}$ and $\psi^\sigma \in \mathcal{L}C_0$. Then $(X^{\underline{\sigma}})^\sigma \vdash_{\mathcal{L}C_0} (\varphi^{\underline{\sigma}})^\sigma$. But $X \subseteq \mathcal{L}$ and $\varphi \in \mathcal{L}$, so $X^{\underline{\sigma}} = X$ and $\varphi^{\underline{\sigma}} = \varphi$. Therefore, $X^\sigma \vdash_{\mathcal{L}C_0} \varphi^\sigma$, so Proposition 5.2.16 gives $X^\sigma \vdash_{\mathcal{L}C} \varphi^\sigma$.

Now suppose $X^\sigma \vdash_{\mathcal{L}C} \varphi^\sigma$. Note that $\psi \in \mathcal{L}$ implies $\psi^\sigma \in \mathcal{L}C'$. Hence, Proposition 5.2.16 implies $X^\sigma \vdash_{\mathcal{L}C'} \varphi^\sigma$. By Theorem 5.2.11, we may choose countable $X_0 \subseteq X$ and $C_0 \subseteq C'$ so that $X_0^\sigma \vdash_{\mathcal{L}C_0} \varphi^\sigma$. Let $V_0 = \{x \in V' \mid c_x \in C_0\}$ and note that V_0 is countable. By Proposition 5.2.8, we may assume $\text{bnd } X_0^\sigma \cup \{\varphi^\sigma\}$ is disjoint from both V_0 and $\text{free } X_0^\sigma \cup \{\varphi^\sigma\}$. We may then use Proposition 5.2.12 to ensure $\text{free } X_0^\sigma \cup \{\varphi^\sigma\}$ is also disjoint from V_0 . Now define $\underline{\sigma}$ as above. If $c_x \in C_0$, then $c_x^{\underline{\sigma}} = x \in V_0$. Therefore, $\underline{\sigma}$ avoids $\text{var } X_0^\sigma \cup \{\varphi^\sigma\}$. Proposition 5.2.15 now gives $(X_0^\sigma)^\sigma \vdash_{\mathcal{L}} (\varphi^\sigma)^\sigma$. Since $\varphi \in \mathcal{L}$ and $\varphi^\sigma \in \mathcal{L}C_0$, we have $(\varphi^\sigma)^\sigma = \varphi$. Similarly, $(X_0^\sigma)^\sigma = X_0$. Therefore, $X_0 \vdash_{\mathcal{L}} \varphi$. \square

5.2.5 Tautologies and consistency

(S:pred-axioms) **Proposition 5.2.19.** *For any $x \in \text{Var}$, we have $\vdash \exists x x = x$.*

Proof. Let $\varphi(x) = (x = x)$. By (i) and (vii)', we have $\forall x \neg \varphi(x) \vdash \neg \varphi(x)$. By (ix) and (ii), we have $\forall x \neg \varphi(x) \vdash \varphi(x)$. Hence, (v) implies $\forall x \neg \varphi(x) \vdash \exists x \varphi(x)$. Since $\exists x \varphi(x) = \neg \forall x \neg \varphi(x)$, we have $\neg \forall x \neg \varphi(x) \vdash \exists x \varphi(x)$, by (i). Therefore, the result follows from (vi). \square

Let $\exists_1 = \exists x_0 x_0 = x_0$. Informally, \exists_1 says that at least one object exists. Justified by Proposition 5.2.19, we define *verum* and *falsum*, respectively, by $\top = \exists_1$ and $\perp = \neg \top$.

More generally, for integers $n > 1$, we define

$$\exists_n = \exists x_0 \cdots x_{n-1} \bigwedge_{i < j < n} x_i \neq x_j.$$

Informally, \exists_n asserts that there are at least n objects. Let $\exists_{=n} = \exists_n \wedge \neg \exists_{n+1}$, which asserts that there are exactly n objects. We also define $\exists_\infty = \bigwedge_{n=1}^\infty \exists_n$, which say that infinitely many objects exist. This last sentence is the only one that is not in \mathcal{L}_{fin} .

A set $X \subseteq \mathcal{L}$ is *inconsistent* if $X \vdash \varphi$ for all $\varphi \in \mathcal{L}$; it is otherwise *consistent*. Note that X is inconsistent if and only if $X \vdash \perp$. Theorem 3.1.13 is easily seen to hold also for \mathcal{L} .

A formula $\varphi \in \mathcal{L}$ is a *tautology* if $\vdash \varphi$; it is a *contradiction* if $\{\varphi\}$ is inconsistent. Note that φ is a tautology if and only if $\neg \varphi$ is a contradiction, and vice versa. As in Proposition 3.1.14, we have $X \vdash \varphi$ if and only if there exists a countable $X_0 \subseteq X$ such that $\vdash \bigwedge X_0 \rightarrow \varphi$.

5.2.6 Deductive and inductive theories

(S:pred-theories) A set $T \subseteq \mathcal{L}^0$ is a (*deductive*) *theory* if $T \vdash \varphi$ implies $\varphi \in T$ for all $\varphi \in \mathcal{L}^0$. The intersection of any family of theories is again a theory. Also, \mathcal{L}^0 itself is a theory. Hence, if $X \subseteq \mathcal{L}^0$, then we may define *the (deductive) theory generated by X* , denoted by $T(X)$ or T_X , as the smallest theory having X as a subset. Note that $T(X) = \{\varphi \in \mathcal{L}^0 \mid X \vdash \varphi\}$.

We adopt the same notation for theories that we did in our propositional calculus. The smallest theory is the set of tautological sentences, which we denote by *Taut*, or *Taut* $_{\mathcal{L}}$. Note that unlike the propositional case, *Taut* is not the set of tautologies. A tautology φ is in *Taut* is and only if $\varphi \in \mathcal{L}^0$. The largest theory is \mathcal{L}^0 . A theory T is inconsistent if and only if $T = \mathcal{L}^0$. The definition of logical equivalence, and its associated notation, are all the same as in the propositional calculus. Note that Lemma 3.1.22 holds also in \mathcal{L}^0 .

The construction of inductive derivability in predicate languages follows exactly as it does in the propositional case. Let $\mathcal{L}^{\text{IS}} = \mathfrak{P} \mathcal{L}^0 \times \mathcal{L}^0 \times [0, 1]$ denote the set of *inductive statements* in \mathcal{L} . All of the results in Sections 3.2–3.5 depend only the fact that $\vdash_{\mathcal{F}}$ satisfies (i)–(vi) of Definition 3.1.3. Since Theorem 5.2.17 shows that $\vdash_{\mathcal{L}}$ restricted to \mathcal{L}^0 also satisfies (i)–(vi), it follows that all of those

results hold in the predicate case as well, with \mathcal{F} replaced by \mathcal{L}^0 . We adopt all of the notation and terminology of Sections 3.2–3.5 to define, in \mathcal{L}^{IS} , inductive theories, inductive conditions, inductive derivability, and all their associated notions.

5.2.7 Karp's calculus

(S:Karp-calc) Karp's completeness theorem [16, Theorem 11.4.1], which we present later in Theorem 5.3.13, will be essential for us. Karp, however, defines her deductive calculus in a different way. We present here Karp's calculus, and show that it is equivalent to the calculus of natural deduction that we defined earlier. We follow the presentation of her calculus that is given in [18, Chapter 4].

Let Λ^- be the smallest subset of \mathcal{L} such that if $x, y \in \text{Var}$, $t \in \mathcal{T}$, $\varphi, \psi \in \mathcal{L}$, and $\Phi \subseteq \mathcal{L}$ is countable with $\varphi \in \Phi$, then the following formulas are in Λ^- :

$$(\Lambda 1) \quad (\varphi \rightarrow \psi \rightarrow \zeta) \rightarrow (\varphi \rightarrow \psi) \rightarrow \varphi \rightarrow \zeta$$

$$(\Lambda 2) \quad (\varphi \rightarrow \neg\psi) \rightarrow \psi \rightarrow \neg\varphi$$

$$(\Lambda 3) \quad \bigwedge \Phi \rightarrow \varphi$$

$$(\Lambda 4) \quad \forall x\varphi \rightarrow \varphi(t/x) \text{ when } t \text{ is free for } x \text{ in } \varphi$$

$$(\Lambda 5) \quad x = x$$

$$(\Lambda 6) \quad x = y \rightarrow y = x$$

$$(\Lambda 7) \quad \varphi \wedge t = x \rightarrow \varphi(t/x) \text{ when } t \text{ is free for } x \text{ in } \varphi$$

The set of *axioms*, or *logical theorems*, denoted by $\Lambda = \Lambda_{\mathcal{L}}$, is the smallest subset of \mathcal{L} such that

$$(I) \quad \Lambda^- \subseteq \Lambda,$$

$$(II) \quad \text{if } \psi, \psi \rightarrow \varphi \in \Lambda, \text{ then } \varphi \in \Lambda,$$

$$(III) \quad \text{if } \psi \rightarrow \varphi \in \Lambda \text{ and } x \notin \text{free } \psi, \text{ then } \psi \rightarrow \forall x\varphi \in \Lambda,$$

$$(IV) \quad \text{If } \Phi \subseteq \mathcal{L} \text{ is countable and } \psi \rightarrow \theta \in \Lambda \text{ for all } \theta \in \Phi, \text{ then } \psi \rightarrow \bigwedge \Phi \in \Lambda.$$

A *proof* of $\varphi \in \mathcal{L}$ from $X \subseteq \mathcal{L}$ is an $(\alpha + 1)$ -sequence of formulas, $\langle \varphi_\beta \mid \beta \leq \alpha \rangle$, where α is a countable ordinal, $\varphi_\alpha = \varphi$, and for each $\beta \leq \alpha$, either $\varphi_\beta \in X \cup \Lambda$, or there exist $i, j < \beta$ such that $\varphi_i = (\varphi_j \rightarrow \varphi_\beta)$, or there exists $\Phi \subseteq \{\varphi_\xi \mid \xi < \beta\}$ such that $\varphi_\beta = \bigwedge \Phi$. Note that if $\langle \varphi_\beta \mid \beta \leq \alpha \rangle$ is a proof of φ_α from X , then for any $\beta < \alpha$, it follows that $\langle \varphi_\xi \mid \xi \leq \beta \rangle$ is a proof of φ_β from X . For $\varphi \in \mathcal{L}$ and $X \subseteq \mathcal{L}$, define $X \vdash \varphi$ to mean there is a proof of φ from X . Note that by (II) above and Proposition 5.2.21 below, $\varphi \in \Lambda$ if and only if $\vdash \varphi$.

(L:add-ante) **Lemma 5.2.20.** *If $\varphi \in \Lambda$, then $\psi \rightarrow \varphi \in \Lambda$ for all $\psi \in \mathcal{L}$.*

Proof. Let $\varphi \in \Lambda$ and $\psi \in \mathcal{L}$. By $(\Lambda 3)$, we have $\neg\neg\psi \wedge \neg\varphi \rightarrow \neg\varphi \in \Lambda$. Thus, $(\Lambda 2)$ and (II) give $\varphi \rightarrow \neg(\neg\neg\psi \wedge \neg\varphi) \in \Lambda$. But unwinding our shorthand shows that $\neg(\neg\neg\psi \wedge \neg\varphi) = \neg\psi \vee \varphi = \psi \rightarrow \varphi$. Therefore, $\varphi \rightarrow \psi \rightarrow \varphi \in \Lambda$. A final application of (II) yields $\psi \rightarrow \varphi \in \Lambda$. \square

$\langle P:\text{conj-axioms} \rangle$ **Proposition 5.2.21.** *The set Λ satisfies*

(IV)' If $\Phi \subseteq \Lambda$ is countable, then $\bigwedge \Phi \in \Lambda$.

Proof. Fix $\theta_0 \in \Phi$. By Lemma 5.2.20, $\theta_0 \in \Lambda$ and $\theta_0 \rightarrow \theta \in \Lambda$ for all $\theta \in \Phi$. Hence, (IV) implies $\theta_0 \rightarrow \bigwedge \Phi \in \Lambda$, and therefore, using (II), we have $\bigwedge \Phi \in \Lambda$. \square

$\langle R:\text{prop-taut} \rangle$ **Remark 5.2.22.** Let \mathcal{F} be a propositional language with infinitely many propositional variables. Given a function $\tau : PV \rightarrow \mathcal{L}$, we extend it to $\tau : \mathcal{F} \rightarrow \mathcal{L}$ recursively by $(\neg\varphi)^\tau = \neg\varphi^\tau$ and $(\bigwedge \Phi)^\tau = \bigwedge_{\varphi \in \Phi} \varphi^\tau$. If $\varphi \in \mathcal{F}_{\text{fin}}$ is a tautology, then we call φ^τ an *instance of a finitary propositional tautology*. By Remark 3.1.15, every such φ^τ can be derived using $(\Lambda 1)$ – $(\Lambda 3)$, (I), (II), and (IV)'. Hence, every instance of a finitary propositional tautology is an axiom.

$\langle P:\text{Hilbert-ded-thm} \rangle$ **Proposition 5.2.23.** *Let $\varphi, \psi \in \mathcal{L}$. Then $\psi \vdash \varphi$ if and only if $\vdash \psi \rightarrow \varphi$.*

Proof. Suppose $\psi \vdash \varphi$ and let $\langle \varphi_\beta \mid \beta \leq \alpha \rangle$ be a proof of φ from ψ . If $\alpha = 0$, then $\varphi = \psi$ or $\varphi \in \Lambda$. If $\varphi = \psi$, then $\vdash \varphi \rightarrow \varphi$ by Remark 5.2.22. If $\varphi \in \Lambda$, then $\vdash \psi \rightarrow \varphi$ by Lemma 5.2.20.

Now assume $\alpha > 0$ and the result is true whenever there is a proof of length less than α . As above, if $\varphi \in \{\psi\} \cup \Lambda$, then $\vdash \psi \rightarrow \varphi$. Suppose $\varphi_i = \varphi_j \rightarrow \varphi$ for some $i, j < \alpha$. Then $\vdash \psi \rightarrow \varphi_j$ and $\vdash \psi \rightarrow \varphi_j \rightarrow \varphi$. By $(\Lambda 1)$,

$$\vdash (\psi \rightarrow \varphi_j \rightarrow \varphi) \rightarrow (\psi \rightarrow \varphi_j) \rightarrow \psi \rightarrow \varphi.$$

With two applications of modus ponens, we obtain $\vdash \psi \rightarrow \varphi$. Finally, suppose $\varphi = \bigwedge \Phi$, where $\Phi \subseteq \{\varphi_\xi \mid \xi < \beta\}$. Then $\vdash \psi \rightarrow \theta$ for all $\theta \in \Phi$. From (IV), it follows that $\vdash \psi \rightarrow \varphi$.

Conversely, suppose $\vdash \psi \rightarrow \varphi$. Then $\psi \vdash \psi \rightarrow \varphi$ and $\psi \vdash \psi$. Applying modus ponens gives $\psi \vdash \varphi$. \square

The proof of Proposition 3.1.16 carries through so that it also holds in this setting.

$\langle T:\text{pred-Hilbert=nat} \rangle$ **Theorem 5.2.24.** *Let $X \subseteq \mathcal{L}$ and $\varphi \in \mathcal{L}$. Then $X \vdash \varphi$ if and only if $X \vdash \varphi$.*

Proof. We first prove that $X \vdash \varphi$ implies $X \vdash \varphi$. It suffices to show that (1)–(3) in Proposition 3.1.16 hold when \vdash is replaced by \vdash . The proof of Theorem 3.1.17 carries through in this case, leaving us only to show that $\vdash \varphi$ whenever $\varphi \in \Lambda$. We prove this by induction using (I)–(IV). Our base case is (I), in which we prove that $\vdash \varphi$ for all $\varphi \in \Lambda^-$.

Axioms of the form $(\Lambda 1)$ and $(\Lambda 2)$ are covered by Remark 5.2.22. Axioms of the form $(\Lambda 3)$ – $(\Lambda 5)$ are covered by (iii), (vii), and (ix), respectively. For $(\Lambda 6)$,

let $s = x$, $t = y$, and $\varphi = (x = z)$. By (x), we have $x = y, x = z \vdash y = z$, which gives

$$\vdash x = z \rightarrow x = y \rightarrow y = z.$$

From here, (viii) gives

$$\vdash \forall z(x = z \rightarrow x = y \rightarrow y = z),$$

so that by (vii),

$$\vdash (x = z \rightarrow x = y \rightarrow y = z)(x/z) = (x = x \rightarrow x = y \rightarrow y = x).$$

Hence, $x = x \vdash x = y \rightarrow y = x$. Combined with $\vdash x = x$, we get ($\Lambda 6$).

Finally, for ($\Lambda 7$), suppose t is free for x in φ . By (x), we have $x = t, \varphi \vdash \varphi(t/x)$. Reversing the roles of x and y in ($\Lambda 6$) gives $\vdash y = x \rightarrow x = y$. From (viii) and (vii), it follows that $\vdash t = x \rightarrow x = t$. Putting these together, we obtain $\vdash t = x \wedge \varphi \rightarrow \varphi(t/x)$, and this completes our base case.

For inductive step (II), suppose $\psi, \psi \rightarrow \varphi \in \Lambda$ and $\vdash \psi, \psi \rightarrow \varphi$. Then we directly have $\vdash \varphi$, and we are done. Step (IV) is equally straightforward. For step (III), suppose $\psi \rightarrow \varphi \in \Lambda$ with $x \notin \text{free } \psi$ and $\vdash \psi \rightarrow \varphi$. Then $\psi \vdash \varphi$, so that (viii) implies $\psi \vdash \forall x\varphi$, which gives $\vdash \psi \rightarrow \forall x\varphi$.

To prove that $X \vdash \varphi$ implies $X \sim \varphi$, it suffices to show that (i)–(x) hold when \vdash is replaced by \sim . The fact that (i)–(vi) hold for \sim follows exactly as in the propositional case. We get (vii) from ($\Lambda 4$). For (viii), suppose $x \notin \text{free } X$ and $X \sim \varphi$. Fix a proof of φ from X and let X_0 be the set of $\theta \in X$ that appear in the proof. Since proofs have countable lengths, X_0 is countable. Also, $X_0 \sim \varphi$. By ($\Lambda 3$), we have $\bigwedge X_0 \sim \varphi$. Lemma 5.2.23 implies $\sim \bigwedge X_0 \rightarrow \varphi$, so that (III) gives $\sim \bigwedge X_0 \rightarrow \forall x\varphi$. Since $X \sim \bigwedge X_0$, an application of modus ponens yields $X \sim \forall x\varphi$.

From ($\Lambda 5$), (viii), and (vii), we obtain (ix). For (x), suppose s and t are free for x in φ and $X \sim s = t, \varphi(s/x)$. By ($\Lambda 7$), (viii), and (vii), we have $\sim \varphi(s/x) \wedge t = s \rightarrow \varphi(t/x)$. Hence, $t = s, \varphi(s/x) \sim \varphi(t/x)$. In the same way, but using ($\Lambda 6$), we get $s = t \sim t = s$. Combining these gives $X \sim \varphi(t/x)$. \square

5.3 Predicate models

(S:pred-models)

In this section, we present the semantics of both deductive and inductive predicate logic. We define satisfiability and consequence, and prove σ -compactness, soundness, and completeness. As with the predicate calculus, the bulk of our work will be in the deductive case. The inductive case will require very little modification from its presentation in Chapter 4.

5.3.1 Strict satisfiability

Let L be an extralogical signature and \mathcal{L} the set of formulas built from L . We will use the two phrases, “ L -structure” and “ \mathcal{L} -structure,” interchangeably. Let

$\omega = (A, L^\omega)$ be an \mathcal{L} -structure. For ground terms t , we define $t^\omega \in A$ recursively by $(ft_1 \cdots t_n)^\omega = f^\omega(t_1^\omega, \dots, t_n^\omega)$.

An *assignment* v into A is a function $v : \mathbf{Var} \rightarrow A$. We can extend v to a function $v_\omega : \mathcal{T} \rightarrow A$ by $v_\omega(c) = c^\omega$ and $v_\omega(ft_1 \cdots t_n) = f^\omega v_\omega(t_1) \cdots v_\omega(t_n)$. The extended function v_ω is called an *assignment into* ω . When there is no risk of confusion, we will omit the subscript and also write v for the extended function v_ω .

Note that if t is a ground term, then $v(t) = t^\omega$. More generally, if $t = t(x_1, \dots, x_n)$, then $v(t)$ depends only on $v(x_1), \dots, v(x_n)$. In this case, we will write $t^\omega[\vec{a}]$, for $\vec{a} \in A^n$, as shorthand for $v(t)$, where v is any assignment satisfying $v(x_i) = a_i$.

Given an assignment v , if $x \in \mathbf{Var}$ and $a \in A$, then we define a new assignment v_x^a by $v_x^a(x) = a$ and $v_x^a(y) = v(y)$ for $y \neq x$.

(D:strict-sat) **Definition 5.3.1.** Let ω be a structure and v an assignment into ω . For $\varphi \in \mathcal{L}$, we define $\omega \models \varphi[v]$ recursively as follows:

- (i) $\omega \models (s = t)[v]$ if and only if $v(s) = v(t)$,
- (ii) $\omega \models (rt_1 \cdots t_n)[v]$ if and only if $r^\omega v(t_1) \cdots v(t_n)$,
- (iii) $\omega \models (\neg\varphi)[v]$ if and only if $\omega \not\models \varphi[v]$,
- (iv) $\omega \models (\bigwedge \Phi)[v]$ if and only if $\omega \models \varphi[v]$ for all $\varphi \in \Phi$, and
- (v) $\omega \models (\forall x\varphi)[v]$ if and only if $\omega \models \varphi[v_x^a]$ for all $a \in A$.

If $\omega \models \varphi[v]$, we say that ω *strictly satisfies* φ *with* v . Note that if $\varphi \in \mathcal{L}_{\text{fin}}$, then \models is the usual notion of satisfiability from first-order logic.

If $\varphi = \varphi(x_1, \dots, x_n)$ and v and v' are assignments that agree on x_1, \dots, x_n , then $\omega \models \varphi[v]$ if and only if $\omega \models \varphi[v']$. In this case, we will write $\omega \models \varphi[\vec{a}]$, where $\vec{a} \in A^n$, to mean that $\omega \models \varphi[v]$ for all assignments v satisfying $v(x_i) = a_i$. In particular, if φ is a sentence, then $\omega \models \varphi$ means that $\omega \models \varphi[v]$ for all assignments v . A formula φ is said to be *strictly satisfiable* if $\omega \models \varphi[v]$ for some structure ω and some assignment v .

The proofs of the following three theorems are the same as in first-order logic, except we use Definition 5.3.1(iv). See, for instance, [28, Theorems 2.3.1, 2.3.4, and 2.3.5] for details.

(T:coinc-thm) **Theorem 5.3.2 (Coincidence theorem).** *Let ω and ω' be structures with a common domain. Let v and v' be assignments into ω and ω' , respectively. Let $\varphi \in \mathcal{L}$ and assume*

- (i) $v(x) = v'(x)$ for all $x \in \text{free } \varphi$, and
- (ii) $s^\omega = s^{\omega'}$ for all $s \in \text{sym } \varphi$.

Then $\omega \models \varphi[v]$ if and only if $\omega' \models \varphi[v']$.

(T:invar-thm) **Theorem 5.3.3 (Invariance theorem).** *Let ω and ω' be isomorphic \mathcal{L} -structures and let $g : \omega \rightarrow \omega'$ be an isomorphism. Let v be an assignment into ω and define the assignment v' into ω' by $v'(x) = gv(x)$. Then $\omega \models \varphi[v]$ if and only if $\omega' \models \varphi[v']$, for all $\varphi \in \mathcal{L}$. In particular, if $\varphi = \varphi(x_1, \dots, x_n)$, then $\omega \models \varphi[\vec{a}]$ if and only if $\omega' \models \varphi[g\vec{a}]$ for all $\vec{a} \in A^n$, where A is the domain of ω . Consequently, if $\varphi \in \mathcal{L}^0$ is a sentence, then $\omega \models \varphi$ if and only if $\omega' \models \varphi$.*

If v is an assignment into ω and σ is a substitution, then v^σ is the assignment defined by $v^\sigma(x) = v(x^\sigma)$. By induction on t , we have $v^\sigma(t) = v(t^\sigma)$ for all $t \in \mathcal{T}$.

Theorem 5.3.4 (Substitution theorem). *Let v be an assignment into the structure ω . Let $\varphi \in \mathcal{L}$ and let σ be a substitution. If σ is free for φ , then $\omega \models \varphi^\sigma[v]$ if and only if $\omega \models \varphi[v^\sigma]$. In particular, if t is free for x in φ , then $\omega \models \varphi(t/x)[v]$ if and only if $\omega \models \varphi[v_x^a]$, where $a = v(t)$.*

Remark 5.3.5. The strict satisfiability relation is not σ -compact. For example, let $L = \{c_n \mid n \in \mathbb{N}_0\}$ be a set of distinct constant symbols and let

$$X = \{\forall x \bigvee_{n \in \mathbb{N}_0} x = c_n\} \cup \{x \neq y \mid x, y \in \text{Var}, x \neq y\}.$$

Then every countable subset of X is strictly satisfiable, but X is not satisfiable.

5.3.2 Models and deductive satisfiability

An *inductive \mathcal{L} -model*, or simply a *model*, is a probability space, $\mathcal{P} = (\Omega, \Sigma, \mathbb{P})$, where Ω is a set of \mathcal{L} -structures. An *assignment into \mathcal{P}* is an indexed collection $\mathbf{v} = \langle v_\omega \mid \omega \in \Omega \rangle$, where v_ω is an assignment into ω for each $\omega \in \Omega$. Note that \mathbf{v} does not depend on Σ or \mathbb{P} . We may therefore sometimes call \mathbf{v} an assignment into Ω .

If \mathbf{v} is an assignment into a model $\mathcal{P} = (\Omega, \Sigma, \mathbb{P})$ and $\varphi \in \mathcal{L}$, then we define

$$\varphi[\mathbf{v}]_\Omega = \{\omega \in \Omega \mid \omega \models \varphi[v_\omega]\}.$$

We say that \mathcal{P} *satisfies φ with \mathbf{v}* , denoted by $\mathcal{P} \models \varphi[\mathbf{v}]$, if $\varphi[\mathbf{v}]_\Omega \in \overline{\Sigma}$ and $\overline{\mathbb{P}}\varphi[\mathbf{v}]_\Omega = 1$. A set $X \subseteq \mathcal{L}$ is *satisfiable* if there is a model \mathcal{P} and an assignment \mathbf{v} into \mathcal{P} such that $\mathcal{P} \models \psi[\mathbf{v}]$ for all $\psi \in X$.

If φ is a sentence, then $\varphi[\mathbf{v}]_\Omega$ does not depend on \mathbf{v} . In this case, we simply write φ_Ω , and as in the propositional case, we have

$$\varphi_\Omega = \{\omega \in \Omega \mid \omega \models \varphi\}.$$

We then write $\mathcal{P} \models \varphi$ to mean $\mathcal{P} \models \varphi[\mathbf{v}]$ for all assignments \mathbf{v} , and this holds if and only if $\varphi_\Omega \in \overline{\Sigma}$ and $\overline{\mathbb{P}}\varphi_\Omega = 1$.

(P:pred-pre-cpct) **Proposition 5.3.6.** *Let $X \subseteq \mathcal{L}$.*

- (i) *If X is strictly satisfiable, then X is satisfiable.*
- (ii) *If X is satisfiable and countable, then X is strictly satisfiable.*

Proof. For (i), suppose $\omega \models \psi[v]$ for all $\psi \in X$. Let $\mathcal{P} = (\{\omega\}, \{\emptyset, \{\omega\}\}, \delta_\omega)$ and $\mathbf{v} = \langle v \rangle$. Then $\mathcal{P} \models \varphi[\mathbf{v}]$. For (ii), suppose X is satisfiable and countable. Let $\mathcal{P} = (\Omega, \Sigma, \mathbb{P})$ be a model and \mathbf{v} an assignment into \mathcal{P} such that $\mathcal{P} \models \psi[\mathbf{v}]$ for all $\psi \in X$. Then $\overline{\mathbb{P}} \bigcap_{\psi \in X} \psi[\mathbf{v}]_\Omega = 1$, so we may choose $\omega \in \bigcap_{\psi \in X} \psi[\mathbf{v}]_\Omega$. We then have $\omega \models \psi[v_\omega]$ for all $\psi \in X$, so that ω strictly satisfies X with v_ω . \square

Given a model $\mathcal{P} = (\Omega, \Sigma, \mathbb{P})$, let

$$\Sigma_{\mathcal{L}} = \Sigma \cap \{\varphi[\mathbf{v}]_\Omega \mid \varphi \in \mathcal{L} \text{ and } \mathbf{v} \text{ is an assignment into } \mathcal{P}\},$$

and let $\mathbb{P}_{\mathcal{L}} = \mathbb{P}|_{\Sigma_{\mathcal{L}}}$. Then $\Sigma_{\mathcal{L}}$ is a sub- σ -algebra of Σ , so that $(\Omega, \Sigma_{\mathcal{L}}, \mathbb{P}_{\mathcal{L}})$ is also a model. For any $\varphi \in \mathcal{L}$ and any assignment \mathbf{v} into \mathcal{P} , we have $\varphi[\mathbf{v}]_\Omega \in \Sigma$ if and only if $\varphi[\mathbf{v}]_\Omega \in \Sigma_{\mathcal{L}}$. Hence, from a logical standpoint, every set $A \in \Sigma \setminus \Sigma_{\mathcal{L}}$ is irrelevant.

Let $\mathcal{P} = (\Omega, \Sigma, \mathbb{P})$ and $\mathcal{Q} = (\Omega', \Gamma, \mathbb{Q})$ be models. We say that \mathcal{P} and \mathcal{Q} are *isomorphic (as models)*, denoted by $\mathcal{P} \simeq \mathcal{Q}$, if there exists a measurable function $h : \Omega \rightarrow \Omega'$ such that h induces an isomorphism (as measure spaces) from $(\Omega, \Sigma_{\mathcal{L}}, \mathbb{P}_{\mathcal{L}})$ to $(\Omega', \Gamma_{\mathcal{L}}, \mathbb{Q}_{\mathcal{L}})$, and $\omega \simeq h\omega$ for \mathbb{P} -a.e. $\omega \in \Omega$. In this case, we abuse notation and say that $h : \mathcal{P} \rightarrow \mathcal{Q}$ is a *model isomorphism*.

Let $h : \mathcal{P} \rightarrow \mathcal{Q}$ be a model isomorphism and \mathbf{v} an assignment into \mathcal{P} . An assignment \mathbf{v}' into \mathcal{Q} is called an *image of \mathbf{v} under h* if, for each $\omega \in \Omega$, there is a function $g_\omega : \omega \rightarrow h\omega$ such that, for \mathbb{P} -a.e. $\omega \in \Omega$, the function g_ω is an isomorphism and $v'_{\omega'}(x) = g_\omega v_\omega(x)$ for all $x \in \text{Var}$.

An image of \mathbf{v} under h always exists. We can construct one as follows. For each $\omega' \in \Omega'$, choose $a_{\omega'}$ in the domain of ω' . For each $\omega \in \Omega$, if $\omega \simeq h\omega'$, then let g_ω be an isomorphism from ω to $h\omega'$. Otherwise, let g_ω map everything to $a_{h\omega}$. We then define $v'_{\omega'}(x) = g_\omega v_\omega(x)$.

(L:iso-thm) **Lemma 5.3.7.** *Let h be a model isomorphism from $\mathcal{P} = (\Omega, \Sigma, \mathbb{P})$ to $\mathcal{Q} = (\Omega', \Gamma, \mathbb{Q})$. Let \mathbf{v} be an assignment into \mathcal{P} and let \mathbf{v}' be an image of \mathbf{v} under h . Then $h^{-1}\varphi[\mathbf{v}']_{\Omega'} = \varphi[\mathbf{v}]_\Omega$, \mathbb{P} -a.s. Consequently, for all $\varphi \in \mathcal{L}$, we have $\varphi[\mathbf{v}]_\Omega \in \overline{\Sigma}$ if and only if $\varphi[\mathbf{v}']_{\Omega'} \in \overline{\Gamma}$, and in this case, $\overline{\mathbb{Q}}\varphi[\mathbf{v}']_{\Omega'} = \overline{\mathbb{P}}\varphi[\mathbf{v}]_\Omega$.*

Proof. Let $\omega' = h\omega$. For all $\omega \in \Omega$, we have $\omega \in h^{-1}\varphi[\mathbf{v}']_{\Omega'}$ if and only if $\omega' \in \varphi[\mathbf{v}']_{\Omega'}$, which holds if and only if $\omega' \models \varphi[\mathbf{v}']$. By Theorem 5.3.3, for \mathbb{P} -a.e. ω , this is equivalent to $\omega \models \varphi[v_\omega]$, which holds if and only if $\omega \in \varphi[\mathbf{v}]_\Omega$. Hence, $h^{-1}\varphi[\mathbf{v}']_{\Omega'} = \varphi[\mathbf{v}]_\Omega$ \mathbb{P} -a.e. Since h also induces a measure-space isomorphism from \mathcal{P} to \mathcal{Q} , we have $\overline{\mathbb{Q}} = \overline{\mathbb{P}} \circ h^{-1}$. Therefore, $\varphi[\mathbf{v}]_\Omega \in \overline{\Sigma}$ if and only if $\varphi[\mathbf{v}']_{\Omega'} \in \overline{\Gamma}$, and in this case, $\overline{\mathbb{Q}}\varphi[\mathbf{v}']_{\Omega'} = \overline{\mathbb{P}}\varphi[\mathbf{v}]_\Omega$. \square

(T:ded-iso-thm) **Theorem 5.3.8 (Deductive isomorphism theorem).** *Let \mathcal{P} and \mathcal{Q} be isomorphic \mathcal{L} -models and let $h : \mathcal{P} \rightarrow \mathcal{Q}$ be a model isomorphism. Let \mathbf{v} be an assignment into \mathcal{P} and let \mathbf{v}' be an image of \mathbf{v} under h . Then $\mathcal{P} \models \varphi[\mathbf{v}]$ if and only if $\mathcal{Q} \models \varphi[\mathbf{v}']$, for all $\varphi \in \mathcal{L}$. In particular, if φ is a sentence, then $\mathcal{P} \models \varphi$ if and only if $\mathcal{Q} \models \varphi$.*

Proof. Let $\mathcal{P} = (\Omega, \Sigma, \mathbb{P})$ and $\mathcal{Q} = (\Omega', \Gamma, \mathbb{Q})$ be isomorphic. Let h , \mathbf{v} , and \mathbf{v}' be as in the statement of the theorem. Suppose $\mathcal{P} \models \varphi[\mathbf{v}]$. Then $\overline{\mathbb{P}}\varphi[\mathbf{v}]_\Omega = 1$.

By Lemma 5.3.7, we have $\overline{\mathbb{Q}}\varphi[\mathbf{v}']_{\Omega'} = \overline{\mathbb{P}}\varphi[\mathbf{v}]_{\Omega} = 1$. Therefore, $\mathcal{Q} \models \varphi[\mathbf{v}']$. Reversing the roles of \mathcal{P} and \mathcal{Q} gives the converse. \square

(R:a.s.-sure) **Remark 5.3.9.** Let $\mathcal{P} = (\Omega, \Sigma, \mathbb{P})$ be a model. Let $\Omega^* \in \Sigma$ with $\mathbb{P}\Omega^* = 1$. Let $\Sigma^* = \{A \cap \Omega^* \mid A \in \Sigma\}$ and $\mathbb{P}^* = \mathbb{P}|_{\Sigma^*}$. Then $\mathcal{P}^* = (\Omega^*, \Sigma^*, \mathbb{P}^*)$ is a model. Choose ω_0 in Ω^* and define $h : \Omega \rightarrow \Omega^*$ by $h\omega = \omega$ if $\omega \in \Omega^*$ and $h\omega = \omega_0$ if $\omega \notin \Omega^*$. It is straightforward to verify that h is measurable and induces an isomorphism (as measure spaces) from $(\Omega, \Sigma_{\mathcal{L}}, \mathbb{P}_{\mathcal{L}})$ to $(\Omega^*, \Sigma_{\mathcal{L}}^*, \mathbb{P}_{\mathcal{L}}^*)$. Hence, h is a model isomorphism and $\mathcal{P} \simeq \mathcal{P}^*$. It follows that if a given property is true almost surely in a model \mathcal{P} , then we can find an isomorphic model in which it is true for every structure ω .

5.3.3 Deductive consequence and soundness

We say that $\varphi \in \mathcal{L}$ is a *consequence* of $X \subseteq \mathcal{L}$, or that X *entails* φ , which we denote by $X \models \varphi$, if $\mathcal{P} \models \varphi[\mathbf{v}]$ whenever $\mathcal{P} \models \psi[\mathbf{v}]$ for all $\psi \in X$. Note that if X is not satisfiable, then it is vacuously true that $X \models \varphi$ for all $\varphi \in \mathcal{L}$. In the case $X = \emptyset$, we write $\models \varphi$, which means that $\mathcal{P} \models \varphi[\mathbf{v}]$ for all \mathcal{P} and \mathbf{v} . If X and φ are sentences, then $X \models \varphi$ if and only if $\mathcal{P} \models X$ implies $\mathcal{P} \models \varphi$.

(R:strict-conseq) **Remark 5.3.10.** If $X \models \varphi$ and $\omega \models \psi[v]$ for all $\psi \in X$, then $\omega \models \varphi[v]$. To see this, simply apply the above definition to $\mathcal{P} = (\{\omega\}, \{\emptyset, \{\omega\}\}, \delta_{\omega})$ and $\mathbf{v} = \langle v \rangle$. In particular, if X and φ are sentences and $\omega \models X$, then $\omega \models \varphi$.

(P:taut-struct) **Proposition 5.3.11.** *Let $\varphi \in \mathcal{L}$. Then $\models \varphi$ if and only if, for all structures ω and all assignments v_{ω} into ω , we have $\omega \models \varphi[v_{\omega}]$.*

Proof. Suppose $\models \varphi$. Then $\mathcal{P} \models \varphi[\mathbf{v}]$ for all \mathcal{P} and \mathbf{v} . Let ω be a structure and v_{ω} an assignment into ω . Define $\mathcal{P} = (\{\omega\}, \{\emptyset, \{\omega\}\}, \delta_{\omega})$ and $\mathbf{v} = \langle v_{\omega} \rangle$. Then \mathcal{P} is a model and \mathbf{v} is an assignment into \mathcal{P} . By hypothesis, $\mathcal{P} \models \varphi[\mathbf{v}]$, which means $\varphi[\mathbf{v}]_{\Omega} = \{\omega\}$. That is, $\omega \models \varphi[v_{\omega}]$.

Conversely, suppose $\omega \models \varphi[v_{\omega}]$ for all structures ω and all assignments v_{ω} into ω . Let $\mathcal{P} = (\Omega, \Sigma, \mathbb{P})$ be a model and \mathbf{v} an assignment into \mathcal{P} . Then $\varphi[\mathbf{v}]_{\Omega} = \Omega$, so $\overline{\mathbb{P}}\varphi[\mathbf{v}]_{\Omega} = 1$. Therefore, $\mathcal{P} \models \varphi[\mathbf{v}]$. \square

(P:free-elim-conseq) **Proposition 5.3.12.** *Let $X \subseteq \mathcal{L}$ and $\varphi \in \mathcal{L}$, and let σ be a free eliminator. Then $X \models_{\mathcal{L}} \varphi$ if and only if $X^{\sigma} \models_{\mathcal{L}C} \varphi^{\sigma}$.*

Proof. Suppose $X \models_{\mathcal{L}} \varphi$. Let $\mathcal{Q} = (\Omega', \Gamma, \mathbb{Q})$ be an $\mathcal{L}C$ -model and \mathbf{v}' an assignment into \mathcal{Q} . Assume $\mathcal{Q} \models \psi^{\sigma}[\mathbf{v}']$ for all $\psi \in X$. For each $\omega' \in \Omega'$, let ω be its \mathcal{L} -reduct. Let $\Omega = \{\omega \mid \omega' \in \Omega'\}$ and let $\mathcal{P} = (\Omega, \Sigma, \mathbb{P})$ be the measure space image of \mathcal{Q} under the function $\omega' \mapsto \omega$. Define an assignment \mathbf{v} into \mathcal{P} by $v_{\omega}(x) = v'_{\omega'}(x)$ if $x^{\sigma} = x$, and $v_{\omega}(x) = c_x^{\omega'}$ if $x^{\sigma} = c_x$. By term induction, $v_{\omega}(t) = v'_{\omega'}(t^{\sigma})$ for all $t \in \mathcal{T}_{\mathcal{L}}$. Then, by formula induction, we obtain $\omega \models \psi[v_{\omega}]$ if and only if $\omega' \models \psi^{\sigma}[v'_{\omega'}]$ for all $\psi \in \mathcal{L}$. Hence, if h denotes the function $\omega' \mapsto \omega$, then $h^{-1}\psi[\mathbf{v}]_{\Omega} = \psi^{\sigma}[\mathbf{v}']_{\Omega'}$, which gives $\overline{\mathbb{P}}\psi[\mathbf{v}]_{\Omega} = \overline{\mathbb{Q}}\psi^{\sigma}[\mathbf{v}']_{\Omega'}$ for all $\psi \in \mathcal{L}$. Therefore, $\mathcal{P} \models \psi[\mathbf{v}]$ for all $\psi \in X$. By hypothesis, this implies $\mathcal{P} \models \varphi[\mathbf{v}]$, which is equivalent to $\mathcal{Q} \models \varphi^{\sigma}[\mathbf{v}']$, and we have $X^{\sigma} \models_{\mathcal{L}C} \varphi^{\sigma}$.

For the converse, suppose $X^\sigma \models_{\mathcal{LC}} \varphi^\sigma$. Let $\mathcal{P} = (\Omega, \Sigma, \mathbb{P})$ be an \mathcal{L} -model and \mathbf{v} an assignment into \mathcal{P} . Assume $\mathcal{P} \models \psi[\mathbf{v}]$ for all $\psi \in X$. For each $\omega \in \Omega$, define the \mathcal{LC} -structure ω' by $\mathbf{s}^{\omega'} = \mathbf{s}^\omega$ if $\mathbf{s} \in L$, and $c_x^{\omega'} = v_\omega(x)$. Define an assignment \mathbf{v}' into \mathcal{P} by $v'_{\omega'}(x) = v_\omega(x)$ for all $x \in \text{Var}$. Then again by term and formula induction, we have $\omega \models \psi[v_\omega]$ if and only if $\omega' \models \psi^\sigma[v'_{\omega'}]$ for all $\psi \in \mathcal{L}$, which as above yields $\mathcal{P} \models \varphi[\mathbf{v}]$. Therefore, $X \models_{\mathcal{L}} \varphi$. \square

(T:pred-Karp-sent) **Theorem 5.3.13 (Karp's completeness theorem).** *Let $\varphi \in \mathcal{L}^0$ be a sentence. Then $\vdash \varphi$ if and only if $\models \varphi$.*

Proof. Karp's completeness theorem first appears in [16, Theorem 11.4.1]. The version we are citing is [18, Theorem 4.3]. There it is shown that if $\varphi \in \mathcal{L}^0$ is a sentence, then $\omega \models \varphi$ for all structures ω if and only if $\varphi \in \Lambda'$, where Λ' is a set of logical axioms described in [18]. We claim that Λ' is the same as Λ , the set of axioms defined in Section 5.2.7. Since $\varphi \in \Lambda$ if and only if $\vdash \varphi$, our statement of Karp's theorem then follows from Theorem 5.2.24 and Proposition 5.3.11.

Keisler's Λ' differs from Λ in only one way. To describe it, we recursively define the shorthand $\sim\varphi$ as follows:

$$\begin{aligned} \sim\varphi &= \neg\varphi \text{ if } \varphi \text{ is prime,} \\ \sim\neg\varphi &= \varphi, \\ \sim\bigwedge\Phi &= \bigvee_{\theta \in \Phi} \neg\theta = \neg\bigwedge_{\theta \in \Phi} \neg\neg\theta, \text{ and} \\ \sim\forall x\varphi &= \exists x\neg\varphi = \neg\forall x\neg\neg\varphi. \end{aligned}$$

Keisler's Λ' includes everything in Λ , as well as all formulas of the form

$$(\Lambda 8) \quad \neg\varphi \leftrightarrow \sim\varphi$$

To check that $\Lambda' = \Lambda$, we must verify that these formulas are already in Λ . We can break this down according to whether φ is prime, $\varphi = \neg\psi$, $\varphi = \bigwedge\Phi$, or $\varphi = \forall x\psi$. Doing this, applying the definition of \sim , and using Theorem 5.2.24, we must check that

$$\begin{aligned} \vdash \neg\varphi &\leftrightarrow \neg\varphi, \\ \vdash \neg\neg\psi &\leftrightarrow \psi, \\ \vdash \neg\bigwedge\Phi &\leftrightarrow \neg\bigwedge_{\theta \in \Phi} \neg\neg\theta, \text{ and} \\ \vdash \neg\forall x\psi &\leftrightarrow \neg\forall x\neg\neg\psi. \end{aligned}$$

The first two are propositional tautologies. The third follows from $\vdash \theta \leftrightarrow \neg\neg\theta$ and Definition 3.1.3(iii),(iv). The fourth follows from $\vdash \psi \leftrightarrow \neg\neg\psi$ and Definition 5.2.5(vii)',(viii). \square

(T:pred-Karp-form) **Theorem 5.3.14 (Karp's theorem for formulas).** *For any formula $\varphi \in \mathcal{L}$, we have $\vdash \varphi$ if and only if $\models \varphi$.*

Proof. Let σ be a full free eliminator, so that φ^σ is a sentence. By Propositions 5.2.18 and 5.3.12, we have $\vdash_{\mathcal{L}} \varphi$ if and only if $\vdash_{\mathcal{L}C} \varphi^\sigma$, and $\models_{\mathcal{L}} \varphi$ if and only if $\models_{\mathcal{L}C} \varphi^\sigma$. Theorem 5.3.13 gives $\vdash_{\mathcal{L}C} \varphi^\sigma$ if and only if $\models_{\mathcal{L}C} \varphi^\sigma$. \square

As in the propositional case, Karp's completeness theorem allows us to prove the result that was described in Remark 5.2.7.

(P:pred-fin-vs-infin)

Proposition 5.3.15. *Let $X \subseteq \mathcal{L}_{\text{fin}}$ and $\varphi \in \mathcal{L}_{\text{fin}}$. If $X \vdash \varphi$, then $X \vdash_{\text{fin}} \varphi$.*

Proof. Let $X \subseteq \mathcal{L}_{\text{fin}}$ and $\varphi \in \mathcal{L}_{\text{fin}}$. Suppose $X \vdash \varphi$. The well-known completeness theorem from first-order logic states that $X \vdash_{\text{fin}} \varphi$ if and only if, for all structures ω and all assignments v into ω , we have $\omega \models \varphi[v]$ whenever $\omega \models \psi[v]$ for all $\psi \in X$. (See, for instance, [28, Theorem 3.2.7]).

Let ω be a structure and v an assignment into ω . Assume that $\omega \models \psi[v]$ for all $\psi \in X$. Choose countable $X_0 \subseteq X$ such that $\vdash \bigwedge X_0 \rightarrow \varphi$. By Theorem 5.3.14, we have $\models \bigwedge X_0 \rightarrow \varphi$. Hence, $\omega \models (\bigwedge X_0 \rightarrow \varphi)[v]$, which means $\omega \models (\bigwedge X_0)[v]$ implies $\omega \models \varphi[v]$. Since $\omega \models \psi[v]$ for all $\psi \in X_0$, it follows that $\omega \models (\bigwedge X_0)[v]$. Therefore, $\omega \models \varphi[v]$. \square

(T:pred-soundness)

Theorem 5.3.16 (Deductive soundness). *Let $X \subseteq \mathcal{L}$ and $\varphi \in \mathcal{L}$. If $X \vdash \varphi$, then $X \models \varphi$.*

Proof. Suppose $X \vdash \varphi$. Let $\mathcal{P} = (\Omega, \Sigma, \mathbb{P})$ be a model and \mathbf{v} an assignment into \mathcal{P} such that $\mathcal{P} \models \psi[\mathbf{v}]$ for all $\psi \in X$. By Theorem 5.2.11, we may choose countable $X_0 \subseteq X$ with $X_0 \vdash \varphi$. Hence, $\vdash \zeta \rightarrow \varphi$, where $\zeta = \bigwedge X_0$. By Theorem 5.3.14, we have $\models \zeta \rightarrow \varphi$, so that $\mathcal{P} \models (\zeta \rightarrow \varphi)[\mathbf{v}]$. That is, $\mathbb{P}(\zeta \rightarrow \varphi)[\mathbf{v}]_\Omega = 1$. But $\mathbb{P}\psi[\mathbf{v}]_\Omega = 1$ for all $\psi \in X$ and $\zeta[\mathbf{v}]_\Omega = \bigcap_{\psi \in X_0} \psi[\mathbf{v}]_\Omega$. Hence, $\mathbb{P}\zeta[\mathbf{v}]_\Omega = 1$, so that $\zeta[\mathbf{v}]_\Omega^c$ is a null set. Since $(\zeta \rightarrow \varphi)[\mathbf{v}]_\Omega = \zeta[\mathbf{v}]_\Omega^c \cup \varphi[\mathbf{v}]_\Omega$, we have $\mathbb{P}\varphi[\mathbf{v}]_\Omega = \mathbb{P}(\zeta \rightarrow \varphi)[\mathbf{v}]_\Omega = 1$. Therefore, $\mathcal{P} \models \varphi[\mathbf{v}]$. Since \mathcal{P} was arbitrary, this shows that $X \models \varphi$. \square

(C:pred-soundness)

Corollary 5.3.17. *If $X \subseteq \mathcal{L}$ is satisfiable, then X is consistent. If X is countable and consistent, then X is strictly satisfiable.*

Proof. Let $X \subseteq \mathcal{L}$. Suppose X is inconsistent. Then $X \vdash \perp$. By Theorem 5.3.16, we have $X \models \perp$. But $\perp_\Omega = \emptyset$, so $\mathcal{P} \not\models \perp[\mathbf{v}]$ for all \mathcal{P} and \mathbf{v} . Hence, X is not satisfiable. For the second part, suppose X is countable and not strictly satisfiable. Let v be an assignment into a structure ω . Then $\omega \not\models (\bigwedge X)[v]$, which implies $\omega \models (\neg \bigwedge X)[v]$. Since ω and v were arbitrary, we have $\models \neg \bigwedge X$. Theorem 5.3.14 then implies $\vdash \neg \bigwedge X$. Thus, $X \vdash \bigwedge X, \neg \bigwedge X$, so that X is inconsistent. \square

5.3.4 Deductive completeness

According to Theorem 5.3.14, we have that φ is a tautology if and only if $\omega \models \varphi[v]$ for all ω and v . Hence, in any model \mathcal{P} , we have $\varphi \vdash \psi$ implies $\varphi[\mathbf{v}]_\Omega \subseteq \psi[\mathbf{v}]_\Omega$ and $\varphi \equiv \psi$ implies $\varphi[\mathbf{v}]_\Omega = \psi[\mathbf{v}]_\Omega$, for any assignment \mathbf{v} into \mathcal{P} .

In Remark 4.1.14, we saw that in the propositional case, we could obtain a converse to the above if we took Ω to be the set of all strict models. That converse was essential to our proof of both deductive and inductive completeness. In the predicate case, we cannot do this, since the collection of all structures is not a set. Instead, we will use the set of structures defined in the proof of the following proposition.

(P:all-structures) **Proposition 5.3.18.** *There exists a set of structures Ω and an assignment \mathbf{v} into Ω such that $\varphi[\mathbf{v}]_\Omega \subseteq \psi[\mathbf{v}]_\Omega$ implies $\varphi \vdash \psi$, and $\varphi[\mathbf{v}]_\Omega = \psi[\mathbf{v}]_\Omega$ implies $\varphi \equiv \psi$. In particular, if φ and ψ are sentences, then $\varphi_\Omega \subseteq \psi_\Omega$ if and only if $\varphi \vdash \psi$, and $\varphi_\Omega = \psi_\Omega$ if and only if $\varphi \equiv \psi$.*

Proof. Let S be the set of all countable, consistent subsets of \mathcal{L} . By Corollary 5.3.17, for each $X \in S$, we may choose a structure $\omega = \omega_X$ and an assignment v_ω into ω such that $\omega \models \zeta[v_\omega]$ for all $\zeta \in X$. Let $\Omega = \{\omega_X \mid X \in S\}$ and let $\mathbf{v} = \langle v_\omega \mid \omega \in \Omega \rangle$.

For the first implication, let $\varphi, \psi \in \mathcal{L}$ and assume $\varphi \not\vdash \psi$. Then Theorem 3.1.13 implies $X = \{\varphi, \neg\psi\}$ is consistent, so that $X \in S$. Hence, with $\omega = \omega_X \in \Omega$ and v defined as above, we have $\omega \models \varphi[v_\omega]$ and $\omega \not\models (\neg\psi)[v_\omega]$. The latter implies $\omega \not\models \psi[v_\omega]$. Thus, $\omega \in \varphi_\Omega[\mathbf{v}]$ and $\omega \notin \psi_\Omega[\mathbf{v}]$, so that $\varphi_\Omega[\mathbf{v}] \not\subseteq \psi_\Omega[\mathbf{v}]$. Reversing the roles of φ and ψ gives the second implication. \square

(T:pred-compactness) **Theorem 5.3.19 (σ -compactness).** *A set $X \subseteq \mathcal{L}$ is satisfiable if and only if every countable subset of X is satisfiable.*

Proof. The only if part is trivial. Suppose every countable subset of X is satisfiable. Assume X is inconsistent. Then $X \vdash \perp$. By Theorem 5.2.11, there exists countable $X_0 \subseteq X$ such that $X_0 \vdash \perp$, implying that X_0 is inconsistent. By Corollary 5.3.17, we have that X_0 is not satisfiable, a contradiction. Hence, X is consistent.

Let Ω and \mathbf{v} be as in Proposition 5.3.18. Let

$$\Sigma = \{\varphi[\mathbf{v}]_\Omega \mid X \vdash \varphi \text{ or } X \vdash \neg\varphi\}.$$

Then Σ is a σ -algebra. If $A \in \Sigma$, choose φ such that $A = \varphi_\Omega[\mathbf{v}]$. Since X is consistent, we cannot have both $X \vdash \varphi$ and $X \vdash \neg\varphi$. We may therefore define $\mathbb{P}A = 1$ if $X \vdash \varphi$ and 0 otherwise. If $A = \varphi_\Omega[\mathbf{v}] = \psi_\Omega[\mathbf{v}]$, then $\varphi \equiv \psi$, by Proposition 5.3.18. Hence, \mathbb{P} is well-defined.

Since X is consistent, $X \not\vdash \perp$. Thus, $\mathbb{P}\emptyset = \mathbb{P}\perp_\Omega = 0$. Conversely, $X \vdash \top$, so $\mathbb{P}\Omega = \mathbb{P}\top_\Omega = 1$.

Now let $\{A_n\}_{n \in \mathbb{N}} \subseteq \Sigma$ be pairwise disjoint, and define $A = \bigcup_n A_n$. For each n , choose φ_n such that $A_n = \varphi_n[\mathbf{v}]_\Omega$, and define $\varphi = \bigvee_n \varphi_n$. Note that $A = \varphi[\mathbf{v}]_\Omega$. Suppose $m \neq n$. Since

$$(\varphi_m \wedge \varphi_n)[\mathbf{v}]_\Omega = A_m \cap A_n = \emptyset = \perp_\Omega,$$

we have $\varphi_m \wedge \varphi_n \equiv \perp$, implying that $X \not\vdash \varphi_m \wedge \varphi_n$. Therefore, either $X \not\vdash \varphi_m$ or $X \not\vdash \varphi_n$. This implies that there is at most one $n \in \mathbb{N}$ with $\mathbb{P}A_n = 1$. Hence,

$\sum_n \mathbb{P} A_n \in \{0, 1\}$ and

$$\begin{aligned} \sum_n \mathbb{P} A_n = 1 & \text{ iff there exists } n \text{ such that } \mathbb{P} A_n = 1 \\ & \text{ iff there exists } n \text{ such that } X \vdash \varphi_n \\ & \text{ iff } X \vdash \varphi \\ & \text{ iff } \mathbb{P} \varphi[\mathbf{v}]_\Omega = \mathbb{P} A = 1, \end{aligned}$$

showing that \mathbb{P} is countably additive. Thus, \mathbb{P} is a measure on (Ω, Σ) with $\mathbb{P} \Omega = 1$, and so $\mathcal{P} = (\Omega, \Sigma, \mathbb{P})$ is a model.

Now let $\varphi \in X$ be arbitrary. Then $X \vdash \varphi$, so that $\varphi[\mathbf{v}]_\Omega \in \Sigma$, and $\mathbb{P} \varphi[\mathbf{v}]_\Omega = 1$. This shows that $\mathcal{P} \models \varphi[\mathbf{v}]$ for all $\varphi \in X$, and X is satisfiable. \square

$\langle \text{C:pred-compactness} \rangle$ **Corollary 5.3.20.** *A set $X \subseteq \mathcal{L}$ is satisfiable if and only if X is consistent.*

Proof. The only if part is Corollary 5.3.17. Suppose X is not satisfiable. By Theorem 5.3.19, there exists a countable subset $X_0 \subseteq X$ that is not satisfiable. By Proposition 5.3.6, the set X_0 is not strictly satisfiable. Hence, by Corollary 5.3.17, the set X_0 is inconsistent, which implies that X is inconsistent. \square

$\langle \text{T:pred-completeness} \rangle$ **Theorem 5.3.21 (Deductive completeness).** *For $X \subseteq \mathcal{L}$ and $\varphi \in \mathcal{L}$, we have $X \models \varphi$ if and only if $X \vdash \varphi$.*

Proof. The if part is Theorem 5.3.16. Suppose $X \not\models \varphi$. Then $X \cup \{\neg\varphi\}$ is consistent, by Theorem 3.1.13. Thus, $X \cup \{\neg\varphi\}$ is satisfiable, by Corollary 5.3.20. Choose \mathcal{P} and \mathbf{v} such that $\mathcal{P} \models \psi[\mathbf{v}]$ for all $\psi \in X \cup \{\neg\varphi\}$. Then $\mathcal{P} \models \psi[\mathbf{v}]$ for all $\psi \in X$, but $\mathcal{P} \not\models \varphi[\mathbf{v}]$. Thus, $X \not\models \varphi$. \square

5.3.5 Peano arithmetic

$\langle \text{S:PA} \rangle$ As an example of deductive predicate logic, we present the theory of Peano arithmetic in the infinitary setting. Let \mathcal{L} be a language that contains a constant symbol $\underline{0}$, a unary function symbol \mathbf{S} , and binary operation symbols $\{+, \cdot\}$. In the language \mathcal{L} , for each $n \in \mathbb{N}$, we use the shorthand $\underline{n} = \mathbf{S} \cdots \mathbf{S} \underline{0}$, where \mathbf{S} is repeated n times.

Define the formulas

$$\begin{aligned} \varphi_1 : \forall x \mathbf{S}x \neq \underline{0} & & \varphi_2 : \forall xy (\mathbf{S}x = \mathbf{S}y \rightarrow x = y) \\ \varphi_3 : \forall x x + \underline{0} = x & & \varphi_4 : \forall xy x + \mathbf{S}y = \mathbf{S}(x + y) \\ \varphi_5 : \forall x x \cdot \underline{0} = \underline{0} & & \varphi_6 : \forall xy x \cdot \mathbf{S}y = x \cdot y + x \end{aligned}$$

For definiteness, we may assume $x = \mathbf{x}_0$ and $y = \mathbf{x}_1$ in the above, so that this is a finite collection of sentences, rather than a family of formulas indexed by $x, y \in \text{Var}$. Note that each $\varphi_i \in \mathcal{L}_{\text{fin}}^0$. If $\varphi = \varphi(x, \vec{y}) \in \mathcal{L}$, define

$$\text{IS}(\varphi) : \forall \vec{y} (\varphi(\underline{0}/x) \wedge \forall x (\varphi \rightarrow \varphi(\mathbf{S}x/x)) \rightarrow \forall x \varphi)$$

Let $\text{IS} = \{\text{IS}(\varphi) \mid \varphi(x, \vec{y}) \in \mathcal{L}\} \subseteq \mathcal{L}^0$ and $\text{IS}_{\text{fin}} = \text{IS} \cap \mathcal{L}_{\text{fin}}^0$. Since $\text{IS}(\varphi)$ has finite length if and only if φ has finite length, we have $\text{IS}_{\text{fin}} = \{\text{IS}(\varphi) \mid \varphi(x, \vec{y}) \in \mathcal{L}_{\text{fin}}\}$.

In first-order logic, $\Lambda_-^{\text{PA}} = \{\varphi_1, \dots, \varphi_6\} \cup \text{IS}_{\text{fin}}$ are the usual axioms of Peano arithmetic. The set IS_{fin} is called the *axiom schema of induction*. We let $\text{PA}_- = T(\Lambda_-^{\text{PA}})$ and $\text{PA}_{\text{fin}} = \text{PA}_- \cap \mathcal{L}_{\text{fin}}^0$. By Proposition 5.3.15, we have $\text{PA}_{\text{fin}} = \{\varphi \in \mathcal{L}_{\text{fin}}^0 \mid \Lambda_-^{\text{PA}} \vdash_{\text{fin}} \varphi\}$. In other words, PA_{fin} is exactly first-order Peano arithmetic.

We also define $\Lambda^{\text{PA}} = \{\varphi_1, \dots, \varphi_6\} \cup \text{IS}$. This differs from Λ_-^{PA} only in the fact that we are allowed to perform induction on infinitary formulas. Let $\text{PA} = T(\Lambda^{\text{PA}})$. Then $\Lambda_-^{\text{PA}} \subseteq \Lambda^{\text{PA}}$ and $\text{PA}_{\text{fin}} \subseteq \text{PA}_- \subseteq \text{PA}$.

Let \mathcal{N} be the standard structure of arithmetic. That is, $\mathcal{N} = (\mathbb{N}_0, 0, \mathbf{S}, +, \cdot)$, where \mathbf{S} is the function $n \mapsto n + 1$. As usual, we will have to rely on context to know whether $\mathbf{S}, +, \cdot$ are referring to objects in the standard structure, or to symbols in the signature of \mathcal{L} . Since $\mathcal{N} \models \Lambda^{\text{PA}}$, we have $\mathcal{N} \models \text{PA}$ by Remark 5.3.10. It is well-known that there are *nonstandard structures of finitary Peano arithmetic*. That is, there exist structures ω such that $\omega \models \text{PA}_{\text{fin}}$ but $\omega \not\models \mathcal{N}$. As it turns out, the analogous statement is still true for PA_- as we see below in Proposition 5.3.22. On the other hand, Proposition 5.3.23 shows that it is not true for PA . In other words, PA completely characterizes the standard structure of arithmetic, meaning that every true statement about arithmetic is provable in PA . Another way to say this, according to completeness, is that if φ is true in the standard structure of arithmetic, then it is true in every model of PA . This is famously not the case for PA_- , thanks to Gödel's first incompleteness theorem (see [28, Theorem 6.5.1]).

$\langle \text{P:models-of-PA-} \rangle$ **Proposition 5.3.22.** *Let $\mathcal{P} = (\Omega, \Sigma, \mathbb{P})$ be a model. Then $\mathcal{P} \models \text{PA}_-$ if and only if $\omega \models \text{PA}_{\text{fin}}$ for \mathbb{P} -a.e. $\omega \in \Omega$. Consequently, $\text{PA}_- \vdash \varphi$ if and only if $\omega \models \Lambda_-^{\text{PA}}$ implies $\omega \models \varphi[v]$ for all ω and all assignments v into ω .*

Proof. Suppose $\mathcal{P} \models \text{PA}_-$. Then $\mathcal{P} \models \Lambda_-^{\text{PA}}$. Since IS_{fin} is countable, so is Λ_-^{PA} . Hence, $\overline{\mathbb{P}}\Omega^* = 1$, where $\Omega^* = \bigcap_{\varphi \in \Lambda_-^{\text{PA}}} \varphi\Omega$. For every $\omega \in \Omega^*$, we have $\omega \models \Lambda_-^{\text{PA}}$, which implies $\omega \models \text{PA}_{\text{fin}}$.

Conversely, suppose $\omega \models \text{PA}_{\text{fin}}$ for \mathbb{P} -a.e. $\omega \in \Omega$. Since $\Lambda_-^{\text{PA}} \subseteq \mathcal{L}_{\text{fin}}^0$, we have $\Lambda_-^{\text{PA}} \subseteq \text{PA}_{\text{fin}}$. Hence, $\omega \models \Lambda_-^{\text{PA}}$ for \mathbb{P} -a.e. $\omega \in \Omega$. This implies that $\Omega^* = \Omega$, \mathbb{P} -a.e. It follows that $\Omega^* \in \overline{\Sigma}$ and $\overline{\mathbb{P}}\Omega^* = 1$. Therefore, $\mathcal{P} \models \Lambda_-^{\text{PA}}$, which gives $\mathcal{P} \models \text{PA}_-$.

For the second claim, the only if direction follows from Theorem 5.3.21 and Remark 5.3.10. For the if direction, suppose $\omega \models \varphi[v]$ for all ω and all assignments v into ω . Let $\mathcal{P} = (\Omega, \Sigma, \mathbb{P}) \models \text{PA}_-$ and let $\mathbf{v} = \langle v_\omega \rangle$ be an assignment into \mathcal{P} . By the above, $\omega \models \Lambda_-^{\text{PA}}$ for \mathbb{P} -a.e. $\omega \in \Omega$. By hypothesis, $\omega \models \varphi[v_\omega]$ for \mathbb{P} -a.e. $\omega \in \Omega$. Hence, $\overline{\mathbb{P}}\varphi[\mathbf{v}]_\Omega = 1$, so that $\mathcal{P} \models \varphi[\mathbf{v}]$. By Theorem 5.3.21, this gives $\text{PA}_- \vdash \varphi$. \square

$\langle \text{P:models-of-PA} \rangle$ **Proposition 5.3.23.** *Let $\mathcal{P} = (\Omega, \Sigma, \mathbb{P})$ be a model. Then $\mathcal{P} \models \text{PA}_\infty$ if and only if $\omega \simeq \mathcal{N}$ for \mathbb{P} -a.e. $\omega \in \Omega$. Consequently, for all $\varphi \in \mathcal{L}^0$, if $\mathcal{N} \models \varphi$, then $\text{PA}_\infty \vdash \varphi$.*

Proof. For the first claim, the if direction follows from the fact that $\omega \simeq \mathcal{N}$ implies $\omega \models \text{PA}_\infty$. For the only if direction, define the formula $\varphi(x) =$

$(\bigvee_{n \in \mathbb{N}_0} x = \underline{n})$. Suppose $\mathcal{P} \models \text{PA}_\infty$. Then $\mathcal{P} \models \Lambda^{\text{PA}} \cup \{\text{IS}(\varphi)\}$, which is countable. Therefore, $\omega \models \Lambda^{\text{PA}} \cup \{\text{IS}(\varphi)\}$ for \mathbb{P} -a.e. $\omega \in \Omega$. Choose any such ω . Note that $n \mapsto \underline{n}^\omega$ is an embedding of \mathcal{N} into ω .

Clearly, $\omega \models \varphi(\underline{0}/x)$. By the definition of \underline{n} , if a is in the domain of ω and $n \in \mathbb{N}_0$, then $\omega \models (x = \underline{n})[a]$ implies $\omega \models (\text{S}x = \underline{n+1})[a]$. Hence, $\omega \models \varphi[a]$ implies $\omega \models \varphi(\text{S}x/x)[a]$. Since a was arbitrary, we have $\omega \models \forall x(\varphi \rightarrow \varphi(\text{S}x/x))$. It therefore follows that $\omega \models \forall x\varphi$. Hence, the map $n \mapsto \underline{n}^\omega$ is surjective, and so it is an isomorphism from \mathcal{N} to ω .

Finally, let $\varphi \in \mathcal{L}^0$ and suppose $\mathcal{N} \models \varphi$. Let $\mathcal{P} = (\Omega, \Sigma, \mathbb{P})$ be a model with $\mathcal{P} \models \text{PA}_\infty$. By the above result and Theorem 5.3.3, we have $\omega \models \varphi$ for \mathbb{P} -a.e. ω . Hence, $\overline{\mathbb{P}}\varphi_\Omega = 1$, so that $\mathcal{P} \models \varphi$. By Theorem 5.3.21, this gives $\text{PA}_\infty \vdash \varphi$. \square

5.3.6 Inductive consequence and completeness

If \mathcal{P} is a model, we define $\text{Th } \mathcal{P} = \{\varphi \in \mathcal{L}^0 \mid \mathcal{P} \models \varphi\}$. The proof of Proposition 4.1.12 is valid here, and shows that $\text{Th } \mathcal{P}$ is a consistent deductive theory. For $(X, \varphi, p) \in \mathcal{L}^{\text{IS}}$, we say that \mathcal{P} *satisfies* (X, φ, p) , denoted by $\mathcal{P} \models (X, \varphi, p)$ if $X \equiv Y \cup \{\psi\}$ for some $Y \subseteq \text{Th } \mathcal{P}$ and some $\psi \in \mathcal{L}^0$ with $\overline{\mathbb{P}}\varphi_\Omega \cap \psi_\Omega / \overline{\mathbb{P}}\psi_\Omega = p$.

As with the inductive calculus, the results in Section 4.2 depend only on deductive completeness and the fact that $\vdash_{\mathcal{F}}$ satisfies (i)–(vi) of Definition 3.1.3. Hence, all of the proofs in that section go through in the predicate case, with \mathcal{F} replaced by \mathcal{L}^0 , “strict model” replaced by “structure,” and \mathbf{B}^{PV} replaced by the set Ω in Proposition 5.3.18. We adopt all of the notation and terminology of Section 4.2 to define inductive consequence in \mathcal{L}^{IS} , extend it to inductive conditions, and establish completeness.

Similarly, all of the results in Sections 4.1.6 and 4.5.1–4.5.4 carry through with the above three replacements. We therefore adopt all of the notation and terminology of those sections to define independence and its related notions.

To all of this, we add the following.

(T:ind-iso-thm) **Theorem 5.3.24 (Inductive isomorphism theorem).** *Let \mathcal{P} and \mathcal{Q} be isomorphic models. Then $\mathcal{P} \models (X, \varphi, p)$ if and only if $\mathcal{Q} \models (X, \varphi, p)$, for all $(X, \varphi, p) \in \mathcal{L}^{\text{IS}}$.*

Proof. Let $\mathcal{P} = (\Omega, \Sigma, \mathbb{P})$ and $\mathcal{Q} = (\Omega', \Gamma, \mathbb{Q})$ be isomorphic and suppose $\mathcal{P} \models (X, \varphi, p)$. Then $X \equiv Y \cup \{\psi\}$, where $\mathcal{P} \models Y$ and $\overline{\mathbb{P}}\varphi_\Omega \cap \psi_\Omega / \overline{\mathbb{P}}\psi_\Omega = p$. By Theorem 5.3.8, we have $\mathcal{Q} \models Y$. Lemma 5.3.7 implies $\overline{\mathbb{Q}}\psi_{\Omega'} = \overline{\mathbb{P}}\psi_\Omega$ and $\overline{\mathbb{Q}}\varphi_{\Omega'} \cap \psi_{\Omega'} = \overline{\mathbb{Q}}(\varphi \wedge \psi)_{\Omega'} = \overline{\mathbb{P}}(\varphi \wedge \psi)_\Omega = \overline{\mathbb{P}}\varphi_\Omega \cap \psi_\Omega$. We therefore have $\overline{\mathbb{Q}}\varphi_{\Omega'} \cap \psi_{\Omega'} / \overline{\mathbb{Q}}\psi_{\Omega'} = p$, so that $\mathcal{Q} \models (X, \varphi, p)$. \square

5.4 Predicate models and random variables

(S:pred-models-RVs) In this section, we discuss the relationship between predicate models and random variables. Here, random variable is meant in the usual sense of measure-theoretic probability theory. That is, a random variable is a measurable function, defined on a probability space, taking values in a measurable space. Our first goal

will be to prove a predicate analogue of Theorem 4.3.1, which states that every probability space is isomorphic to a propositional model. We begin by establishing the connection between propositional and predicate models.

(P:prop-embed) **Proposition 5.4.1.** *Every propositional model is isomorphic to a predicate model. More specifically, let \mathcal{F} be a given propositional language with propositional variables PV , and let \mathcal{P} be a model in \mathcal{F} . Then there exists a predicate language \mathcal{L} and an \mathcal{L} -model \mathcal{Q} such that $\mathcal{P}_{\mathcal{F}}$ and \mathcal{Q} are isomorphic as measure spaces.*

Proof. Let $\mathcal{P} = (\Omega, \Sigma, \mathbb{P})$ be a propositional model in \mathcal{F} . Let $\alpha = |PV|$ and write $PV = \langle \mathbf{r}_\delta \mid \delta < \alpha \rangle$. Let $\{r_\delta \mid \delta < \alpha\}$ be a set of distinct unary relation symbols. Let ρ be a constant symbol, which we call the propositional constant. Define the extralogical signature $L = \{\rho\} \cup \{r_\delta \mid \delta < \alpha\}$, and let \mathcal{L} be the associated predicate language.

Given a strict propositional model $\omega \in \Omega$, we define the \mathcal{L} -structure $\bar{\omega}$ as follows. The domain of $\bar{\omega}$ will be $A = \mathbf{B}^{PV}$. Let $r_\delta^{\bar{\omega}} = \{\nu \in A \mid \nu \models_{\mathcal{F}} \mathbf{r}_\delta\}$, and let $\rho^{\bar{\omega}} = \omega$. Let $\bar{\Omega} = \{\bar{\omega} \mid \omega \in \Omega\}$ and let $\mathcal{Q} = (\bar{\Omega}, \Gamma, \mathbb{Q})$ be the measure space image of $\mathcal{P}_{\mathcal{F}} = (\Omega, \bar{\Sigma}_{\mathcal{F}}, \bar{\mathbb{P}}_{\mathcal{F}})$ under the function h mapping ω to $\bar{\omega}$. Then \mathcal{Q} is an \mathcal{L} -model.

Define $\tau : PV \rightarrow \mathcal{L}$ by $\mathbf{r}_\delta^\tau = r_\delta \rho$. Extend τ recursively to \mathcal{F} by $(\neg\varphi)^\tau = \neg\varphi^\tau$ and $(\bigwedge \Phi)^\tau = \bigwedge_{\varphi \in \Phi} \varphi^\tau$. We then have $\omega \models_{\mathcal{F}} \varphi$ if and only if $\bar{\omega} \models_{\mathcal{L}} \varphi^\tau$, for all $\varphi \in \mathcal{F}$. This is clear by construction when $\varphi \in PV$. It then follows easily by formula induction on φ .

Now let $A \in \bar{\Sigma}_{\mathcal{F}} = \bar{\Sigma} \cap \mathcal{B}^{PV}$. Choose $\varphi \in \mathcal{F}$ such that $A = \varphi_\Omega$. Define $U = \varphi_\Omega^\tau$. Then $\omega \in A$ if and only if $\omega \models_{\mathcal{F}} \varphi$, and $\omega \in h^{-1}U$ if and only if $\bar{\omega} \models_{\mathcal{L}} \varphi^\tau$. Hence, $A = h^{-1}U$, so that $U \in \Gamma$. Since A was arbitrary, this shows that h induces a measure-space isomorphism from $\mathcal{P}_{\mathcal{F}}$ to \mathcal{Q} . \square

5.4.1 Random variables as extralogical symbols

(S:RV-symb) In Theorem 4.3.1, we showed that every probability space is isomorphic to a propositional model. Conversely, every propositional model is a probability space. In this sense, then, measure-theoretic probability theory is exactly the semantics of propositional inductive logic. But this simple observation misses an important point. While the propositional version of inductive logic is capable of representing any probability space, it does not explicitly represent any random variables.

The reason this matters is that measure-theoretic probability theory is more than just probability spaces. The modern practitioner almost always specializes in a particular class of random variables and stochastic processes. For this reason, we define the following. A *measure-theoretic probability model* is a tuple, (S, Γ, ν, X) , where (S, Γ, ν) is a probability space, $X = \langle X_i \mid i \in I \rangle$ is an indexed collection of random variables, and $\Gamma = \sigma(\langle X_i \mid i \in I \rangle)$. That is, for each $i \in I$, there is a measurable space (R_i, Γ_i) such that $X_i : S \rightarrow R_i$ is measurable, and Γ is the smallest σ -algebra on S that contains $\{X_i \in V\}$ for every $i \in I$ and every $V \in \Gamma_i$.

We aim to prove a predicate analogue of Theorem 4.3.1, and show that every measure-theoretic probability model has a natural correspondence to an \mathcal{L} -model, where the logical signature L is directly connected to the random variables X .

We construct the logical signature as follows. Let $R = \bigcup_{i \in I} R_i$. Let $\{\underline{r} \mid r \in R\}$ be a set of distinct constant symbols, and $\{\underline{V} \mid i \in I, V \in \Gamma_i\}$ a set of distinct unary relation symbols. Let

$$L_R = \{\underline{r} \mid r \in R\} \cup \{\underline{V} \mid i \in I, V \in \Gamma_i\},$$

and let \mathcal{L}_R be the associated predicate language. Define the L_R -structure $\mathcal{R} = (R, L^{\mathcal{R}})$ by $\underline{r}^{\mathcal{R}} = r$ and $\underline{V}^{\mathcal{R}} = V$. Let $T_R = \{\varphi \in \mathcal{L}^0 \mid \mathcal{R} \models \varphi\}$. Then T_R is a deductive theory. In \mathcal{L}_R , we write $y \in \underline{V}$ as shorthand for $\underline{V}y$, and $y \notin \underline{V}$ as shorthand for $\neg \underline{V}y$. Let $C = \{\underline{X}_i \mid i \in I\}$ be a set of distinct constant symbols not in L_R , and define $L = L_R C$.

(T:prob-model-iso) **Theorem 5.4.2.** *There exists an \mathcal{L} -model $\mathcal{P} = (\Omega, \Sigma, \mathbb{P})$ with $\mathcal{P} \models T_R$, and a function $h : S \rightarrow \Omega$ mapping $x \in S$ to $\omega \in \Omega$ such that*

- (i) $x \in \{X_i \in V\}$ if and only if $\omega \models \underline{X}_i \in \underline{V}$,
- (ii) each $U \in \Gamma$ can be written as $U = h^{-1}\varphi_\Omega$ for some $\varphi \in \mathcal{L}^0$, and
- (iii) h induces a measure-space isomorphism from (S, Γ, ν) to \mathcal{P} .

Consequently, if $P = \mathbf{Th} \mathcal{P} \upharpoonright_{[T_R, Th \mathcal{P}]}$, then

$$P(\bigwedge_{k=1}^n \underline{X}_{i(k)} \in \underline{V}_k \mid T_R) = \nu \bigcap_{k=1}^n \{X_{i(k)} \in V_k\}, \quad (5.4.1) \quad \boxed{\text{prob-model-iso}}$$

whenever $i(1), \dots, i(n) \in I$ and $V_k \in \Gamma_{i(k)}$.

Proof. For each $x \in S$, define $\omega = \omega^x$ to be the L -expansion of \mathcal{R} given by $\omega^{\underline{X}_i} = X_i(x)$. Let $\Omega = \{\omega^x \mid x \in S\}$ and let $h : S \rightarrow \Omega$ denote the map $x \mapsto \omega^x$. Let $\mathcal{P} = (\Omega, \Sigma, \mathbb{P})$ be the measure space image of (S, Γ, ν) under h . Since $\omega \models T_R$ for all $\omega \in \Omega$, we have $\mathcal{P} \models T_R$. By construction, we have $X_i(x) \in V$ if and only if $\omega^x \models \underline{X}_i \in \underline{V}$, so (i) holds.

For (ii), let

$$\Gamma' = \{U \in \Gamma \mid U = h^{-1}\varphi_\Omega \text{ for some } \varphi \in \mathcal{L}^0\}.$$

Since $\bigcup_n h^{-1}(\varphi_n)_\Omega = h^{-1}(\bigvee_n \varphi_n)_\Omega$ and $\perp_\Omega = \emptyset$, we have that Γ' is a σ -algebra. Let $V \in \Gamma_i$. Since $X_i(x) \in V$ if and only if $\omega^x \models \underline{X}_i \in \underline{V}$, it follows that $\{X_i \in V\} = h^{-1}(\underline{X}_i \in \underline{V})_\Omega$. Therefore, $\{X_i \in V\} \in \Gamma'$. Since $\Gamma = \sigma(\{X_i \mid i \in I\})$, this proves that $\Gamma = \Gamma'$, so (ii) holds.

If $U \in \Gamma$, $\varphi \in \mathcal{L}^0$, and $U = h^{-1}\varphi_\Omega$, then by the construction of \mathcal{P} , we have $\varphi_\Omega \in \Sigma$. Therefore, (ii) implies (iii).

Finally, since h also induces an isomorphism from $(S, \bar{\Gamma}, \bar{\nu})$ to $(\Omega, \bar{\Sigma}, \bar{\mathbb{P}})$, we have $\bar{\mathbb{P}} = \bar{\nu} \circ h^{-1}$. This gives

$$P(\bigwedge_{k=1}^n \underline{X}_{i(k)} \in \underline{V}_k \mid T_R) = \bar{\mathbb{P}} \bigcap_{k=1}^n (\underline{X}_{i(k)} \in \underline{V}_k)_\Omega = \nu \bigcap_{k=1}^n \{X_{i(k)} \in V_k\},$$

which verifies (5.4.1). \square

We excluded the case where $\sigma(\langle X_i \mid i \in I \rangle)$ is a proper subset of Γ . If we wish to treat this case, we can simply add ι , the identity function on S , to our list of random variables. Note, however, that in this case, there are events $U \in \Gamma$ that have nothing to do with any of the random variables X_i . These are analogous to propositional sentences in the sense that they are generic assertions that lack structure. If we add ι to our list of random variables, and $\rho = \iota \in L$ is the constant symbol that represents ι , then ρ is playing the same role as the propositional constant in the proof of Proposition 5.4.1. We see, then, that we are effectively treating every event $U \in \Gamma \setminus \sigma(\langle X_i \mid i \in I \rangle)$ as if it were a propositional variable.

Theorem 5.4.2 shows that every measure-theoretic probability model is an inductive model. In other words, the whole of measure-theoretic probability theory is embedded in the semantics of inductive logic. But Theorem 5.4.2 says more than just this. It exhibits a particular embedding. The function h in Theorem 5.4.2 gives us a logical interpretation for each component of a measure-theoretic probability model. With this interpretation, we have the following correspondences.

<i>Measure Theory</i>	<i>Inductive Logic</i>
outcome	structure
event	sentence
set membership	strict satisfiability
random variable	constant symbol

5.4.2 Extralogical symbols as functions

(S:symb-func) An \mathcal{L} -structure is, in fact, a function whose domain is L . If ω is an \mathcal{L} -structure, then it maps each $s \in L$ to the object s^ω . Hence, if $\mathcal{P} = (\Omega, \Sigma, \mathcal{P})$ is an \mathcal{L} -model, then each structure $\omega \in \Omega$ is a function that maps the symbol \underline{X} to the object \underline{X}^ω . This is exactly the opposite of what we have in measure-theoretic probability theory, where each random variable X is a function that maps the outcome ω to the object $X(\omega)$.

Starting with an \mathcal{L} -model, $\mathcal{P} = (\Omega, \Sigma, \mathbb{P})$, we may wish to reverse the natural direction of the mapping, and think of the extralogical symbols $s \in L$ as functions defined on Ω . We can do that as follows. If A_ω is the domain of $\omega \in \Omega$, then a constant symbol c gives rise to the function $X^c(\omega) = c^\omega$, mapping Ω to $A = \bigcup_{\omega \in \Omega} A_\omega$. An n -ary relation symbol can be viewed as an indexed collection of $\{0, 1\}$ -valued functions, indexed by A^n . Namely, for each $\vec{a} \in A^n$, we have $X_{\vec{a}}^r(\omega) = 1$ if $\vec{a} \in r^\omega$, and 0 otherwise. For an n -ary function symbol f , we can add a so-called ‘‘cemetery point’’ to A . Let ∂ be an object not in A . Then f provides us with an indexed collection of $A \cup \{\partial\}$ -valued functions, indexed by A^n . That is,

$$X_{\vec{a}}^f(\omega) = \begin{cases} f^\omega \vec{a} & \text{if } \vec{a} \in A_\omega^n, \\ \partial & \text{if } \vec{a} \notin A_\omega^n. \end{cases}$$

These functions are, of course, not measurable. In fact, the set A is not even equipped with a σ -algebra, and there may not be a natural σ -algebra on A

that is compatible with Σ . Without measurability, we cannot use the well-established theory of random variables to analyze the functions determined by the extralogical symbols. On the other hand, without the requirement of measurability, we are able to model situations that are not possible with random variables. See, for instance, Example 5.4.8 below.

5.4.3 The relativity of randomness

(S:rel-rand) Let $\mathcal{P} = (\Omega, \Sigma, \mathbb{P})$ be an \mathcal{L} -model and let $\underline{X} \in L$ be a constant symbol. If we try to think of $\omega \mapsto \underline{X}^\omega$ as a kind of non-measurable random variable, then we run into a problem deeper than its non-measurability. The problem we face is that *every* extralogical symbol is a non-measurable random variable. In a measure-theoretic probability model, if we are faced with a probability of the form $\nu\{X > 0\}$, then we can be quite certain that the only thing random is X . But in an inductive model, the analogous expression is $\overline{\mathbb{P}}\{\omega \in \Omega \mid X^\omega >^\omega 0^\omega\}$. Not only can the value of 0 vary with ω , the inequality relation itself can also depend on ω .

This phenomenon can be seen in a very simple example. Let $L = \{\underline{h}, \underline{t}, \underline{X}\}$ be a set of constant symbols and \mathcal{L} the associated predicate language. We think of \underline{h} and \underline{t} as denoting the heads and tails sides of a coin, and \underline{X} the result of flipping the coin. Let $T_0 \subseteq \mathcal{L}^0$ be the deductive theory generated by the sentences, $\underline{h} \neq \underline{t}$ and $\underline{X} = \underline{h} \vee \underline{X} = \underline{t}$. We may think of T_0 as describing our state of knowledge prior to flipping the coin. Namely, the two sides of the coin are distinct, and the coin will not land on its edge.

Let P be the inductive theory generated by

$$P(\underline{X} = \underline{h} \mid T_0) = P(\underline{X} = \underline{t} \mid T_0) = 1/2.$$

This, of course, represents our assumption that the coin is fair.

Intuitively, we imagine that \underline{h} and \underline{t} are fixed, whereas \underline{X} is random. We can satisfy P with a model that matches this intuition. Let $A = \{0, 1\}$. Define the L -structure ω_0 by $\underline{h}^{\omega_0} = 1$, $\underline{t}^{\omega_0} = 0$, and $\underline{X}^{\omega_0} = 0$. Define the L -structure ω_1 by $\underline{h}^{\omega_1} = 1$, $\underline{t}^{\omega_1} = 0$, and $\underline{X}^{\omega_1} = 1$. Let $\Omega = \{\omega_0, \omega_1\}$, $\Sigma = \mathfrak{P}\Omega$, and $\mathbb{P}\{\omega_0\} = \mathbb{P}\{\omega_1\} = 1/2$. Then $\mathcal{P} = (\Omega, \Sigma, \mathbb{P}) \models P$.

Under \mathcal{P} , the symbol \underline{h} corresponds to the function $\omega \mapsto \underline{h}^\omega = 1$, and the symbol \underline{t} corresponds to the function $\omega \mapsto \underline{t}^\omega = 0$. In other words, \underline{h} and \underline{t} are identified with constant functions, and are therefore fixed. On the other hand, \underline{X} corresponds to $\omega \mapsto \underline{X}^\omega$, which is 0 with probability 1/2 and 1 with probability 1/2. Hence, \underline{X} is random.

However, we can also satisfy P with a model that violates this intuition. Let ω_0 be as above. Define ω'_1 by $\underline{h}^{\omega'_1} = 0$, $\underline{t}^{\omega'_1} = 1$, and $\underline{X}^{\omega'_1} = 0$. Let $\Omega' = \{\omega_0, \omega'_1\}$, $\Sigma' = \mathfrak{P}\Omega'$, and $\mathbb{P}'\{\omega_0\} = \mathbb{P}'\{\omega'_1\} = 1/2$. Then $\mathcal{P}' = (\Omega', \Sigma', \mathbb{P}') \models P$. This time, however, \underline{h} and \underline{t} correspond to functions that are 0 or 1 with equal probability, and \underline{X} corresponds to the constant function 0. In this model, it is \underline{h} and \underline{t} that are random, while \underline{X} is fixed.

Since P is satisfied by both \mathcal{P} and \mathcal{P}' , we see that P does not tell us which terms are random and which terms are fixed. In fact, it is not even meaningful

to ask this question in P . The only things in P which can be random (that is, the only things which can be assigned a probability that is not 0 or 1) are sentences. In order to even ask this question, we must fix a model. And in fixing a model, we are adopting, so to speak, a point of view. Which terms are random and which are fixed is relative to that point of view. In \mathcal{P} , we take the point of view that \underline{h} and \underline{t} are fixed, while \underline{X} is random. And in \mathcal{P}' , we take the point of view that \underline{X} is fixed, while \underline{h} and \underline{t} are random. There are models in which all three are random. In this example, however, there are no models in which all three are fixed.

In general, then, whether a term in P is random or fixed depends on our point of view, or to borrow the language of physics, it depends on our frame of reference.

5.4.4 Frames of reference

A *frame of reference* is a method that takes a given \mathcal{L} -model \mathcal{P} and constructs a new \mathcal{L} -model \mathcal{P}' such that $\mathcal{P} \simeq \mathcal{P}'$. Formally, we could define a frame of reference to be a class function that maps each \mathcal{P} in a certain class of \mathcal{L} -models to a set of \mathcal{L} -models that are isomorphic to \mathcal{P} . This level of formalism, however, will not be necessary for our purposes.

Let P be the inductive theory in Section 5.4.3 that models a fair coin flip. Proposition 5.4.6 below gives a method of taking any \mathcal{L} -model \mathcal{P} such that $\mathcal{P} \models P$, and constructing an isomorphic model \mathcal{P}' in which the functions $\omega \mapsto \underline{h}^\omega$ and $\omega \mapsto \underline{t}^\omega$ are constant functions. In other words, there is a frame of reference in which \underline{h} and \underline{t} are fixed, and not random.

We begin by showing there is a frame of reference in which every object is an ordinal. A model $\mathcal{P} = (\Omega, \Sigma, \mathbb{P})$ is said to be an *ordinal model* if, for all structures $\omega \in \Omega$, the domain of ω is an ordinal.

(L:ord-FOR) **Lemma 5.4.3.** *Let $\mathcal{P} = (\Omega, \Sigma, \mathbb{P})$ be a model. For each $\omega \in \Omega$, let A_ω be the domain of ω , and let B_ω be a set with $|B_\omega| = |A_\omega|$. Choose a bijection $g_\omega : A_\omega \rightarrow B_\omega$, and let ω' be the isomorphic image of ω under g_ω . Let $\Omega' = \{\omega' \mid \omega \in \Omega\}$ and let $h : \Omega \rightarrow \Omega'$ denote the function $\omega \mapsto \omega'$. Define $\mathcal{Q} = (\Omega', \Gamma, \mathbb{Q})$ to be the measure space image of \mathcal{P} under h . Then h is a model isomorphism from \mathcal{P} to \mathcal{Q} .*

Proof. To verify that h is an isomorphism from \mathcal{P} to \mathcal{Q} , it suffices to check that h induces an isomorphism as measure spaces from $(\Omega, \Sigma_{\mathcal{L}}, \mathbb{P}_{\mathcal{L}})$ to $(\Omega', \Gamma_{\mathcal{L}}, \mathbb{Q}_{\mathcal{L}})$. For this, it suffices to show that for all $A \in \Sigma_{\mathcal{L}}$, there exists $U \in \Gamma_{\mathcal{L}}$ such that $h^{-1}U = A$.

Let $A \in \Sigma_{\mathcal{L}}$. Choose $\varphi \in \mathcal{L}$ and choose an assignment \mathbf{v} into \mathcal{P} such that $A = \varphi[\mathbf{v}]_\Omega$. Define the assignment \mathbf{v}' into \mathcal{Q} by $v'_{\omega'}(x) = g_\omega v_\omega(x)$, and let $U = \varphi[\mathbf{v}']_{\Omega'}$. It now suffices to show that $U \in \Gamma_{\mathcal{L}}$ and $h^{-1}U = A$. But \mathcal{Q} is the measure space image of \mathcal{P} under h . Hence, if $h^{-1}U = A \in \Sigma_{\mathcal{L}} \subseteq \Sigma$, then $U \in \Gamma$, which implies $U \in \Gamma_{\mathcal{L}}$ by the definition of $\Gamma_{\mathcal{L}}$. Therefore, we need only show that $h^{-1}U = A$.

Note that $\omega \in h^{-1}U = h^{-1}\varphi[v']_{\Omega'}$ if and only if $\omega' \models \varphi[v'_\omega]$. Similarly, $\omega \in A = \varphi[v]_{\Omega}$ if and only if $\omega \models \varphi[v_\omega]$. By Theorem 5.3.3, we have $\omega \models \varphi[v]$ if and only if $\omega' \models \varphi[v']$. Thus, $h^{-1}U = A$, and so h is an isomorphism. \square

(P:ord-FOR) **Proposition 5.4.4 (Ordinal frame of reference).** *Every model is isomorphic to an ordinal model.*

Proof. Let $\mathcal{P} = (\Omega, \Sigma, \mathbb{P})$ be a model. For each $\omega \in \Omega$, let A_ω be the domain of ω . Choose an ordinal α_ω such that $|\alpha_\omega| = |A_\omega|$, and choose a bijection $g_\omega : A_\omega \rightarrow \alpha_\omega$. Define \mathcal{Q} as in Lemma 5.4.3. Then $\mathcal{P} \simeq \mathcal{Q}$ and \mathcal{Q} is an ordinal model. \square

(L:fix-constants) **Lemma 5.4.5.** *Let α be an ordinal and $S \subseteq \alpha$. Let β and γ be ordinals such that $|\beta| = |S|$ and $|\gamma| = |\alpha \setminus S|$, and let $g : S \rightarrow \beta$ be a bijection. Then g can be extended to a bijection $g : \alpha \rightarrow \beta + \gamma$.*

Proof. Let $g : S \rightarrow \beta$ be a bijection. Choose a bijection $h : \alpha \setminus S \rightarrow \gamma$. Note that the function $f : \gamma \rightarrow (\beta + \gamma) \setminus \beta$ given by $f\xi = \beta + \xi$ is a bijection. Therefore, $f \circ h : \alpha \setminus S \rightarrow (\beta + \gamma) \setminus \beta$ is a bijection. Hence, if we define $g\xi = fh\xi$ for $\xi \in \alpha \setminus S$, then $g : \alpha \rightarrow \beta + \gamma$ is a bijection. \square

(P:fix-constants) **Proposition 5.4.6 (Constant frame of reference).** *Let \mathcal{L} be a predicate language with extralogical signature L . Let $C = \{c_0, c_1, \dots\} \subseteq L$ be a countable (possibly finite) set of constant symbols. Let $T \subseteq \mathcal{L}^0$ be a deductive theory. Assume that $T \vdash c_m \neq c_n$ for all $m \neq n$. Then for all models \mathcal{P} such that $\mathcal{P} \models T$, there exists an ordinal model $\mathcal{P}' = (\Omega', \Sigma', \mathbb{P}')$ such that $\mathcal{P} \simeq \mathcal{P}'$ and $c_n^{\omega'} = n$ for every $\omega' \in \Omega'$.*

Proof. Suppose that $\mathcal{P} \models T$. By Proposition 5.4.4, we may assume that \mathcal{P} is an ordinal model. Let $\varphi = (\bigwedge_{m \neq n} c_m \neq c_n)$. We then have $T \vdash \varphi$, so that $\mathcal{P} \models \varphi$. Let $\omega \in \Omega$ and let α_ω denote the domain of ω . We define an ordinal α'_ω and a bijection $g_\omega : \alpha_\omega \rightarrow \alpha'_\omega$ as follows. If $\omega \notin \varphi_\Omega$, then let $\alpha'_\omega = \alpha_\omega$ and let g_ω be the identity. Suppose $\omega \in \varphi_\Omega$. Then $\omega \models \varphi$, which means $c_m^\omega \neq c_n^\omega$ for all $m \neq n$.

Let $\beta = |C|$, so that either $\beta = \{0, 1, \dots, N\}$ or $\beta = \mathbb{N}_0$. Then $C = \{c_n \mid n \in \beta\}$. Define $S = \{c_n^\omega \mid n \in \beta\} \subseteq \alpha_\omega$. Since $c_m^\omega \neq c_n^\omega$ for all $m \neq n$, we have $|\beta| = |S|$ and $n \mapsto c_n^\omega$ is a bijection from β to S . Define $g_\omega c_n^\omega = n$, so that $g_\omega : S \rightarrow \beta$ is a bijection. Choose an ordinal γ such that $|\gamma| = |\alpha_\omega \setminus S|$ and define $\alpha'_\omega = \beta + \gamma$. By Lemma 5.4.5, we may choose an extension of g_ω to α_ω such that $g_\omega : \alpha_\omega \rightarrow \alpha'_\omega$ is a bijection.

Having constructed α'_ω and g_ω , we now define the ordinal model $\mathcal{Q} = (\Omega', \Gamma, \mathbb{Q})$ as in Lemma 5.4.3, so that $\mathcal{P} \simeq \mathcal{Q}$. By Theorem 5.3.8, we have $\mathcal{Q} \models \varphi$. Hence, $\overline{\mathbb{Q}}\varphi_{\Omega'} = 1$. Let $\omega' \in \varphi_{\Omega'}$. Then $\omega' \models \varphi$, which implies $\omega \models \varphi$, since $\omega' \simeq \omega$. Therefore, $\omega \in \varphi_\Omega$, and it follows that $c_n^{\omega'} = g_\omega c_n^\omega = n$. Hence, $c_n^{\omega'} = n$ for \mathbb{Q} -a.e. $\omega' \in \Omega'$. By Remark 5.3.9, we may assume that $c_n^{\omega'} = n$ for every $\omega' \in \Omega'$. \square

5.4.5 The natural frame of reference

Consider an inductive theory P with root T_0 such that $\text{PA}_- \subseteq T_0$. In P , we may have inductive statements of the form $P(\underline{X} > \underline{n} \mid T_0) = p$, where $>$ is the extralogical symbol defined by $\forall xy(x > y \leftrightarrow (\exists z \neq 0) x = y + z)$. In any model $\mathcal{P} = (\Omega, \Sigma, \mathbb{P})$ that satisfies P , we then have $\mathbb{P}(\underline{X} > \underline{n})_\Omega = p$. But

$$(\underline{X} > \underline{n})_\Omega = \{\omega \in \Omega \mid \underline{X}^\omega >^\omega \underline{n}^\omega\}.$$

Hence, the very meaning of \underline{n} and $>$ in the model \mathcal{P} may vary with ω . If we carry with us an intuition that was built around a study of random variables, then this situation is highly counterintuitive. We are not accustomed to thinking of positive integers as random, let alone thinking of $>$ as a random relation. But recall that the randomness or fixedness of these symbols is not an inherent property of the inductive theory P that we started it. It is relative to the model we are considering.

Theorem 5.4.7 below, which is an immediate consequence of Proposition 5.4.6, shows that for any such inductive theory P , there is a frame of reference in which all the constant symbols \underline{n} are fixed. We call this the *natural frame of reference*. If we replace PA_- with PA , then this also fixes $>$. Indeed, in this case, $>$ can be defined explicitly by

$$\forall xy(x > y \leftrightarrow \bigvee_{n > m} (x = \underline{n} \wedge y = \underline{m})).$$

This is because $\text{PA} \vdash \forall x \bigvee_{n \in \mathbb{N}_0} x = \underline{n}$, which was demonstrated in the proof of Proposition 5.3.23.

To state the formal theorem, let \mathcal{L} be a language that contains a unary function symbol \mathbf{S} and constant symbols $\{\underline{n} \mid n \in \mathbb{N}_0\}$. A deductive theory $T \subseteq \mathcal{L}^0$ is said to *contain the counting numbers* if

$$\begin{aligned} T &\vdash \forall x \mathbf{S}x \neq \underline{0}, \\ T &\vdash \forall xy(\mathbf{S}x = \mathbf{S}y \rightarrow x = y), \text{ and} \\ T &\vdash \underline{n} = \mathbf{S} \cdots \mathbf{S}\underline{0}, \text{ for all } n \in \mathbb{N}_0. \end{aligned}$$

In the last condition, the symbol \mathbf{S} is repeated n times.

(T:fix-naturals) **Theorem 5.4.7 (Natural frame of reference).** *Let $T \subseteq \mathcal{L}^0$ be a deductive theory that contains the counting numbers. Then for all models \mathcal{P} such that $\mathcal{P} \models T$, there exists an ordinal model $\mathcal{P}' = (\Omega', \Sigma', \mathbb{P}')$ such that $\mathcal{P} \simeq \mathcal{P}'$ and $\underline{n}^\omega = n$ for every $\omega \in \Omega'$.*

Proof. Since T contains the counting numbers, we have $T \vdash \underline{n} \neq \underline{m}$ for all $n \neq m$. The theorem therefore follows from Proposition 5.4.6. \square

(Exp1:pred-Karp412) **Example 5.4.8.** Let I be an uncountable set. Let

$$L = \{\underline{X}_t \mid t \in I\} \cup \{\underline{n} \mid n \in \mathbb{N}_0\}$$

be a set of constant symbols, and \mathcal{L} the associated predicate language. Define $X \subseteq \mathcal{L}^0$ by

$$X = \{\underline{m} \neq \underline{n} \mid m, n \in \mathbb{N}_0, m \neq n\} \cup \{\bigvee_{n \in \mathbb{N}_0} \underline{X}_t = \underline{n} \mid t \in I\} \\ \cup \{\underline{X}_s \neq \underline{X}_t \mid s, t \in I, s \neq t\}.$$

Let $X_0 \subseteq X$ be countable and choose $I_0 = \{t_0, t_1, t_2, \dots\}$ such that

$$X_0 \subseteq \{\underline{m} \neq \underline{n} \mid m, n \in \mathbb{N}_0, m \neq n\} \cup \{\bigvee_{n \in \mathbb{N}_0} \underline{X}_t = \underline{n} \mid t \in I_0\} \\ \cup \{\underline{X}_s \neq \underline{X}_t \mid s, t \in I_0, s \neq t\}.$$

Define the \mathcal{L} -structure ω with domain $A = \mathbb{N}_0$ by $\underline{n}^\omega = n$, $\underline{X}_t^\omega = n$ if $t = t_n$, and $\underline{X}_t^\omega = 0$ otherwise. Then $\omega \models X_0$, so that X_0 is satisfiable, by Proposition 5.3.6. Since X_0 was arbitrary, Theorem 5.3.19 implies X is satisfiable. Choose a model $\mathcal{P} = (\Omega, \Sigma, \mathbb{P})$ such that $\mathcal{P} \models X$. By Proposition 5.4.6, we may assume that $\underline{n}^\omega = n$ for all $\omega \in \Omega$.

Note that for each $t \in I$, we have $\mathcal{P} \models \bigvee_{n \in \mathbb{N}_0} \underline{X}_t = \underline{n}$. Hence, $\underline{X}_t^\omega \in \mathbb{N}_0$ for \mathbb{P} -a.e. $\omega \in \Omega$. Moreover, if $T_0 = T(X)$ and $P = \mathbf{Th} \mathcal{P} \upharpoonright_{[T_0, T_h \mathcal{P}]}$, then $P(\underline{X}_s \neq \underline{X}_t \mid T_0) = 1$ for all $s \neq t$. This should be contrasted with the observation made in Remark 4.4.7. Namely, there is no \mathbb{N}_0 -valued stochastic process $\langle Y(t) \mid t \in I \rangle$ such that $Y(s) \neq Y(t)$ a.s. for all $s \neq t$.

Chapter 6

Real inductive theories

(Ch:real-ind-ths) By a “real inductive theory,” we mean an inductive theory that makes statements about real numbers. If $P \subseteq \mathcal{L}^{\text{IS}}$ is such a theory with root T_0 , then \mathcal{L} should be capable of making statements about real numbers, and (ideally) T_0 should contain all true statements about real numbers.

One particularly straightforward way to construct such an inductive theory is to follow the approach taken in Section 5.4.1. Namely, we construct a standard structure of the real numbers, which we denote by \mathcal{R} , and we require T_0 to contain all sentences that are strictly satisfied by \mathcal{R} . We can then prove the analogue of Theorem 5.4.2, showing that every collection of real-valued random variables can be represented in a natural way inside a real inductive theory. This is done in Section 6.3.5.

There are several downsides to this approach. The first is that we cannot talk directly about sets of real numbers. We can add them indirectly as relations in our language, as we did in Section 5.4.1. But there is no intrinsic theory of sets in this language. A second, related downside is that we cannot talk about distinguished elements or subsets of the reals without considerable extra effort. In particular, we cannot talk directly about integers and rationals, and their relationships to the reals.

The primary purpose of this chapter is to present a different, more robust approach. Namely, we will create inductive theories whose root T_0 contains all of axiomatic set theory. In this way, not only can we make inductive statements about real numbers, but also about all other objects of modern mathematics.

After discussing definitorial extensions in Section 6.1, the axioms of set theory are presented in Section 6.2. They are the usual axioms of Zermelo-Fraenkel set theory with choice. As with Peano arithmetic in Section 5.3.5, we define multiple theories. In fact, we define four theories: $\text{ZFC}_{\text{fin}} \subseteq \text{ZFC}_- \subseteq \text{ZFC} \subseteq \text{ZFC}_+$. The first of these, ZFC_{fin} , is the usual finitary set theory from first-order logic. The others are extensions to \mathcal{L}^0 . The first extension, ZFC_- , is conservative, in the sense that every sentence in $\text{ZFC}_- \setminus \text{ZFC}_{\text{fin}}$ is purely infinitary. In other words, in ZFC_- , we cannot deduce any new first-order sentences that we could not already deduce in ZFC_{fin} . This is because ZFC_- and ZFC_{fin} have

the same axioms. In particular, even in ZFC_- , we can only use finitary formulas when making use of the axioms of separation and replacement (see Section 6.2.2 and 6.2.3).

The extension ZFC is stronger. There, we allow infinitary formulas in the axiom of separation. In ZFC_+ , we allow infinitary formulas in both separation and replacement. We do not spend time on ZFC_+ beyond defining it. We primarily focus on the theories ZFC_- and ZFC .

In Section 6.3, we construct the set of real numbers in ZFC_- , and build real inductive theories whose roots are required to contain ZFC_- . We then prove the analogue of Theorem 5.4.2, showing that every collection of real-valued random variables can be represented inside such a theory. Using ZFC_- is superior to using the standard structure of the reals. Now, not only can we talk about distinguished sets of real numbers, we can in fact talk about all kinds of sets, backed up by the full power of set theory.

These benefits, however, come at a price. In ZFC_- , although we can define the set of real numbers, we cannot define each individual real number. We can explicitly define each rational number, and we can explicitly define certain individual real numbers, such as π , e , and $\sqrt{2}$. But it is intuitively clear, at least in finitary set theory, that the vast majority of real numbers elude any type of description. As such, our probabilistic statements will only be able to mention the rationals, and a handful of definable reals.

To elaborate on this, consider an inductive statement, $P(\varphi \mid X) = p$. There are two things to notice about this. First, it is not an element of \mathcal{L} . We cannot add a quantifier to the outside of this statement, except in a metatheoretical sense. Second, the formula φ is a sentence. It cannot contain any free variables. Hence, any mention of a real number inside φ must be done through an explicitly defined constant. Therefore, when using ZFC_- , any real number which cannot be explicitly defined in ZFC_- cannot be used inside an inductive statement.

Another downside to using ZFC_- is that it produces a weaker analogue of Theorem 5.4.2. In that theorem, we see a direct and intuitive correspondence between outcomes in the measure-theoretic model, and structures in the inductive model. This connection is lost when we do things in ZFC_- .

It turns out that the right place to work is in ZFC . This is done in Section 6.4. In ZFC , not only can we explicitly define each individual real number, we can also explicitly define each individual Borel set, and each individual measurable function. Hence, any statement we might make in our measure-theoretic model has an explicit counterpart in ZFC . Moreover, we recover the natural correspondence between outcomes and structures.

Additionally, in ZFC , we can construct a frame of reference in which the real numbers, Borel sets, and measurable functions are all almost surely fixed, and not random. In this sense, ZFC is home to the natural intuition of the practicing probabilist, to whom it would never occur to think of such things as varying with ω .

Adopting ZFC , however, involves accepting a new axiom of set theory—or rather, accepting an expanded version of the axiom of separation. In Section 6.2.7, we discuss reasons why this is hardly any more problematic than assuming

that ZFC_{fin} is consistent.

Finally, in Sections 6.5 and 6.6, we illustrate how the major theorems and structures of measure-theoretic probability can be expressed using inductive logic. The examples we cover are the law of large numbers, the central limit theorem, conditional expectation, and the general form of the law of total probability, also known as the tower property of conditional expectation.

6.1 Definitorial extensions

(S: def-ext) We often want to introduce new symbols into our language that are defined in terms of old ones. Sometimes this can be done using shorthand. In that case, the new symbols are not actually part of our language. They are just notational conventions we use to talk about our language. We have seen this already with the symbols \exists and \rightarrow .

Sometimes, however, we want to formally augment our logical signature. For example, in the context of a deductive theory T , suppose we have $T \vdash \exists! x \varphi(x)$. We may wish to introduce a constant symbol to denote the unique object whose existence is being asserted. This is not easily done with shorthand. In this subsection, we go over precisely how this is done, and what effects it has on deductive and inductive derivability.

6.1.1 Defining individual symbols

Relation symbols. Let \mathcal{L} be a predicate language with logical signature L . Let r be an n -ary relation symbol with $r \notin L$, and let $\mathcal{L}[r]$ be the language with signature $L \cup \{r\}$. An *explicit definition of r in \mathcal{L}* is a sentence in $\mathcal{L}[r]^0$ of the form $\theta_r = \forall x(r\vec{x} \leftrightarrow \delta(\vec{x}))$, where $\delta = \delta(x_1, \dots, x_n) \in \mathcal{L}$. The formula δ is called a *defining formula*. We may sometimes denote δ by δ_r , to indicate its relationship to r .

Given $\varphi \in \mathcal{L}[r]$, we define $\varphi^{\text{rd}} \in \mathcal{L}$ as follows. If φ is an equation, then $\varphi^{\text{rd}} = \varphi$, and if $\varphi = r\vec{t}$, then $\varphi^{\text{rd}} = \delta(\vec{t})$. We extend this recursively by $(\neg\varphi)^{\text{rd}} = \neg\varphi^{\text{rd}}$, $(\bigwedge \Phi)^{\text{rd}} = \bigwedge_{\varphi \in \Phi} \varphi^{\text{rd}}$, and $(\forall x\varphi)^{\text{rd}} = \forall x\varphi^{\text{rd}}$. Intuitively, φ is reduced down to φ^{rd} by replacing all occurrences of $r\vec{t}$ by $\delta(\vec{t})$.

Constant symbols. Let c be a constant symbol with $c \notin L$. Let $\mathcal{L}[c]$ be the language with signature $L \cup \{c\}$. An *explicit definition of c in \mathcal{L}* is a sentence in $\mathcal{L}[c]^0$ of the form $\theta_c = \forall y(y = c \leftrightarrow \delta(y))$, where $\delta = \delta(y) \in \mathcal{L}$. The formula δ is called a *defining formula*. We may sometimes denote δ by δ_c , to indicate its relationship to c . Let $\xi_c = \exists! y \delta_c$. Note that $\theta_c \vdash \xi_c$. In general, we will only use the definition θ_c in situations where ξ_c holds. Given $\varphi \in \mathcal{L}[c]$, we choose $z \notin \text{var } \varphi$ and define $\varphi^{\text{rd}} \in \mathcal{L}$ by $\varphi^{\text{rd}} = \exists z(\varphi(z/c) \wedge \delta(z))$.

Function symbols. Finally, let f be an n -ary function symbol with $f \notin L$ and $n \geq 1$. Let $\mathcal{L}[f]$ be the language with signature $L \cup \{f\}$. An *explicit definition of f in \mathcal{L}* is a sentence in $\mathcal{L}[f]^0$ of the form $\theta_f = \forall \vec{x}(y = f\vec{x} \leftrightarrow \delta(\vec{x}, y))$, where

$\delta = \delta(x_1, \dots, x_n, y) \in \mathcal{L}$. The formula δ is called a *defining formula*. We may sometimes denote δ by δ_f , to indicate its relationship to f . Let $\xi_f = \forall \vec{x} \exists! y \delta_f$. Note that $\theta_f \vdash \xi_f$. In general, we will only use the definition θ_f in situations where ξ_f holds.

Given $\varphi \in \mathcal{L}[f]$, we define $\varphi^{\text{rd}} \in \mathcal{L}$ by formula recursion. First suppose φ in prime. Then φ is a string of finite length. Here, we define φ^{rd} exactly as in first-order logic. (See, for example, [28, Section 2.6].) Namely, choose $y \notin \text{var } \varphi$. Find the leftmost occurrence of f in φ , which will be followed by a unique concatenation of terms $\vec{t} = t_1 \cdots t_n$, and let φ' be the prime formula obtained by replacing $f\vec{t}$ with y . Note that $\varphi = \varphi'(f\vec{t}/y)$. We then define $\varphi_1 = \exists y(\varphi' \wedge \delta(\vec{t}, y))$. The resulting formula φ_1 has one fewer occurrence of f than φ . If f still occurs in φ_1 , then repeat the procedure to obtain φ_2 , and so on. Since φ has only finitely many occurrences of f , this procedure will eventually terminate in some φ_m that no longer contains f . We then define $\varphi^{\text{rd}} = \varphi_m$. We extend this definition recursively by $(\neg\varphi)^{\text{rd}} = \neg\varphi^{\text{rd}}$, $(\bigwedge \Phi)^{\text{rd}} = \bigwedge_{\varphi \in \Phi} \varphi^{\text{rd}}$, and $(\forall x\varphi)^{\text{rd}} = \forall x\varphi^{\text{rd}}$.

6.1.2 Defining multiple symbols

(S: def-ext-mult) More generally, let M be a set of extralogical symbols, disjoint from L . Let \mathcal{L}' be the language with signature $L \cup M$. For each $s \in M$, let θ_s be an explicit definition of s in \mathcal{L} , and let $\Theta = \{\theta_s \mid s \in M\}$. Let $\xi_r = \top$ for all relation symbols $r \in M$, and let $\Xi = \{\xi_s \mid s \in M\}$. Note that $\Theta \vdash \Xi$. In general, we will only use the definitions Θ in situations where Ξ holds.

Given $\varphi \in \mathcal{L}'$, we define the *reduced formula*, $\varphi^{\text{rd}} \in \mathcal{L}$, as follows. If φ is prime, then $\text{sym } \varphi$ is finite. We may therefore eliminate the symbols in $\text{sym } \varphi \cap M$ in a stepwise fashion as above. We then extend this recursively by $(\neg\varphi)^{\text{rd}} = \neg\varphi^{\text{rd}}$, $(\bigwedge \Phi)^{\text{rd}} = \bigwedge_{\varphi \in \Phi} \varphi^{\text{rd}}$, and $(\forall x\varphi)^{\text{rd}} = \forall x\varphi^{\text{rd}}$. More generally, for $X \subseteq \mathcal{L}'$, we write $X^{\text{rd}} = \{\varphi^{\text{rd}} \mid \varphi \in X\}$.

6.1.3 Extensions and models

Let $\omega = (A, L^\omega)$ be an \mathcal{L} -structure, and define the \mathcal{L}' -structure, $\omega' = (A, (L')^{\omega'})$ as follows. First, let $s^{\omega'} = s^\omega$ whenever $s \in L$. If $s = r \in M$ is a relation symbol, then we define $r^{\omega'}$ by $r^{\omega'} \vec{a}$ if and only if $\omega \models \delta_r[\vec{a}]$, where $\vec{a} \in A^n$. Next, suppose $s = c \in M$ is a constant symbol. If $\omega \models \xi_c = \exists! y \delta_c$, then there exists a unique $a \in A$ such that $\omega \models \delta_c[a]$. We then define $c^{\omega'} = a$. Otherwise, if $\omega \not\models \xi_c$, then choose $a \in A$ arbitrarily and set $c^{\omega'} = a$. Lastly, suppose $s = f \in M$ is a function symbol. If $\omega \models \xi_f = \forall \vec{x} \exists! y \delta_f$, then for each $\vec{a} \in A^n$, there exists a unique $b \in A$ such that $\omega \models \delta_f[\vec{a}, b]$. We then define $f^{\omega'}(\vec{a}) = b$. Otherwise, if $\omega \not\models \xi_f$, then we define $f^{\omega'}$ arbitrarily. Note that ω' is constructed so that $\omega' \models \theta_s$ whenever $\omega \models \xi_s$. Conversely, note that if ν is an \mathcal{L}' -model and ω is its \mathcal{L} -reduct, then $\omega \models \xi_s$ whenever $\nu \models \theta_s$, and ν and ω' agree on $L \cup \{s \in M \mid \nu \models \theta_s\}$. Finally, given an assignment v_ω into ω , we define the assignment $v_{\omega'}$ into ω' by $v_{\omega'}(x) = v_\omega(x)$ for all $x \in \text{Var}$.

(P:elim-struct) **Proposition 6.1.1.** *Let ω be an \mathcal{L} -structure and $\varphi \in \mathcal{L}'$. Assume $\omega \models \xi_s$ for all $s \in M \cap \text{sym } \varphi$. Then $\omega' \models \varphi[v'_{\omega'}]$ if and only if $\omega \models \varphi^{\text{rd}}[v_\omega]$, for all assignments v_ω .*

Proof. If $\varphi \in \mathcal{L}$, then $\varphi^{\text{rd}} = \varphi$. Hence, by Theorem 5.3.2, the proposition holds for all $\varphi \in \mathcal{L}$.

For $\varphi \in \mathcal{L}'$, we first consider the case, $M \cap \text{sym } \varphi = \{r\}$. In this case, $\omega \models \xi_r$. It follows that $\omega' \models \theta_r = \forall \vec{x}(r\vec{x} \leftrightarrow \delta(\vec{x}))$. We now prove the proposition by induction on φ . Suppose φ is prime. Then either $\varphi \in \mathcal{L}$ or $\varphi = r\vec{t}$. In the former case, we established that the proposition holds. In the latter case, we have $\varphi^{\text{rd}} = \delta(\vec{t})$, so the result follows from $\omega' \models \theta_r$. The inductive steps are straightforward. The cases $M = \{c\}$ and $M = \{f\}$ are similar.

We now consider the case of general M . As above, the result holds if $\text{sym } \varphi \cap M$ contains a single element. It therefore holds whenever $\text{sym } \varphi \cap M$ is finite, by reducing each symbol one at a time. In particular, it holds for each prime φ , since prime formulas are finite strings of symbols. The result then follows by induction on φ , using the recursive definition of strict satisfiability. \square

Let $\mathcal{P} = (\Omega, \Sigma, \mathbb{P})$ be an \mathcal{L} -model. For each $\omega \in \Omega$, define the \mathcal{L}' -structure ω' as above. Let $\Omega' = \{\omega' \mid \omega \in \Omega\}$, let h denote the map $\omega \mapsto \omega'$, and let $\mathcal{P}' = (\Omega', \Gamma, \mathbb{Q})$ be the measure space image of \mathcal{P} under h . Since $\omega' \models \theta_s$ whenever $\omega \models \xi_s$, it follows that $\mathcal{P}' \models \theta_s$ whenever $\mathcal{P} \models \xi_s$. Given an assignment \mathbf{v} into \mathcal{P} , define the assignment \mathbf{v}' into \mathcal{P}' by $v'_{\omega'}(x) = v_\omega(x)$ for all $x \in \text{Var}$.

(L:elim-model) **Lemma 6.1.2.** *Let $\mathcal{Q} = (\Omega^{\mathcal{Q}}, \Gamma^{\mathcal{Q}}, \mathbb{Q}^{\mathcal{Q}})$ be an \mathcal{L}' -model such that $\mathcal{Q} \models \Theta$. For each $\nu \in \Omega^{\mathcal{Q}}$, let $g\nu$ be the L -reduct of ν . Define $\Omega = \{g\nu \mid \nu \in \Omega^{\mathcal{Q}}\}$ and let $\mathcal{P} = (\Omega, \Sigma, \mathbb{P})$ be the measure space image of \mathcal{Q} under g . Given an assignment \mathbf{w} into \mathcal{Q} , define the assignment \mathbf{v} into \mathcal{P} by $v_{g\nu}(x) = w_\nu(x)$ for all $x \in \text{Var}$. Let \mathcal{P}' and \mathbf{v}' as above. Then $\mathcal{P} \models \Xi$ and, for any $\varphi \in \mathcal{L}'$, we have $\mathcal{Q} \models \varphi[\mathbf{w}]$ if and only if $\mathcal{P}' \models \varphi[\mathbf{v}']$.*

Proof. Let $\xi_s \in \Xi$. Since $\xi_s \in \mathcal{L}$, it follows from Theorem 5.3.2 that $g\nu \models \xi_s$ if and only if $\nu \models \xi_s$. Hence, $(\xi_s)_{\Omega^{\mathcal{Q}}} = g^{-1}(\xi_s)_\Omega$, so that $\mathcal{Q} \models \xi_s$ if and only if $\mathcal{P} \models \xi_s$. Since $\Theta \vdash \Xi$ and $\mathcal{Q} \models \Theta$, this gives $\mathcal{P} \models \Xi$.

Now let $M_0 = M \cap \text{sym } \varphi$ and note that M_0 is countable. Let $\Theta_0 = \{\theta_s \mid s \in M_0\}$. Then $\mathcal{Q} \models \bigwedge \Theta_0$, so that $\nu \models \bigwedge \Theta_0$, for a.e. ν . It follows that ν and ω' agree on $L \cup M_0$, for a.e. ν . The result now follows from Theorem 5.3.2. \square

(P:elim-model) **Proposition 6.1.3.** *Let $\mathcal{P} = (\Omega, \Sigma, \mathbb{P})$ be an \mathcal{L} -model such that $\mathcal{P} \models \Xi$, and let $\mathcal{P}' = (\Omega', \Gamma, \mathbb{Q})$ be as above. Let $\varphi \in \mathcal{L}'$ and let \mathbf{v} be an assignment into \mathcal{P} . Then $\varphi[\mathbf{v}']_{\Omega'} \in \bar{\Gamma}$ if and only if $\varphi^{\text{rd}}[\mathbf{v}]_\Omega \in \bar{\Sigma}$, and in this case, $\bar{\mathbb{Q}}\varphi[\mathbf{v}']_{\Omega'} = \bar{\mathbb{P}}\varphi^{\text{rd}}[\mathbf{v}]_\Omega$. In particular, $\mathcal{P}' \models \varphi[\mathbf{v}']$ if and only if $\mathcal{P} \models \varphi^{\text{rd}}[\mathbf{v}]$.*

Proof. Since $\text{sym } \varphi$ is countable and $\mathcal{P} \models \Xi$, we may choose $\Omega^* \in \Sigma$ such that $\mathbb{P}\Omega^* = 1$ and $\omega \models \xi_s$, for all $s \in M \cap \text{sym } \varphi$ and all $\omega \in \Omega^*$. Hence, by Proposition 6.1.1, we have $\omega' \models \varphi[v'_{\omega'}]$ if and only if $\omega \models \varphi^{\text{rd}}[v_\omega]$, for \mathbb{P} -a.e. $\omega \in \Omega$.

Therefore, $h^{-1}\varphi[\mathbf{v}']_{\Omega'} = \varphi^{\text{rd}}[\mathbf{v}]_{\Omega}$ a.s. It follows from the definition of \mathcal{P}' that $\varphi[\mathbf{v}']_{\Omega'} \in \bar{\Gamma}$ if and only if $\varphi^{\text{rd}}[\mathbf{v}]_{\Omega} \in \bar{\Sigma}$, and in this case, $\mathbb{Q}\varphi[\mathbf{v}']_{\Omega'} = \mathbb{P}\varphi^{\text{rd}}[\mathbf{v}]_{\Omega}$. In particular, $\mathbb{Q}\varphi[\mathbf{v}']_{\Omega'} = 1$ if and only if $\mathbb{P}\varphi^{\text{rd}}[\mathbf{v}]_{\Omega} = 1$, so that $\mathcal{P}' \models \varphi[\mathbf{v}']$ if and only $\mathcal{P} \models \varphi^{\text{rd}}[\mathbf{v}]$. \square

6.1.4 Deductive elimination

The following theorem captures the exact relationship between derivability in \mathcal{L}' and derivability in \mathcal{L} . By Remark 5.2.6, since $\mathcal{L} \subseteq \mathcal{L}'$, we are able to simply write \vdash , instead of $\vdash_{\mathcal{L}}$ and $\vdash_{\mathcal{L}'}$.

$\langle \text{T:ded-elim-thm} \rangle$ **Theorem 6.1.4 (Deductive elimination theorem).** *Let $X \subseteq \mathcal{L}'$ and $\varphi \in \mathcal{L}'$. Then $X, \Theta \vdash \varphi$ if and only if $X^{\text{rd}}, \Xi \vdash \varphi^{\text{rd}}$.*

Proof. Assume $X, \Theta \vdash \varphi$. Let $\mathcal{P} = (\Omega, \Sigma, \mathbb{P})$ be an \mathcal{L} -model and let \mathbf{v} be an assignment into \mathcal{P} such that $\mathcal{P} \models \psi[\mathbf{v}]$ for all $\psi \in X^{\text{rd}} \cup \Xi$. Since $\mathcal{P} \models \Xi$, we may define \mathcal{P}' as in Proposition 6.1.3. We then have $\mathcal{P}' \models \zeta[\mathbf{v}']$ if and only $\mathcal{P} \models \zeta^{\text{rd}}[\mathbf{v}]$ for all $\zeta \in \mathcal{L}'$. It follows that $\mathcal{P}' \models \psi[\mathbf{v}']$ for all $\psi \in X$. Since $\mathcal{P}' \models \theta_s$ whenever $\mathcal{P} \models \xi_s$, it also follows that $\mathcal{P}' \models \Theta$. Thus, $\mathcal{P} \models \psi[\mathbf{v}']$ for all $\psi \in X \cup \Theta$. Since $X, \Theta \vdash \varphi$, we conclude that $\mathcal{P}' \models \varphi[\mathbf{v}']$. One more application of Proposition 6.1.3 gives $\mathcal{P} \models \varphi^{\text{rd}}[\mathbf{v}]$. Since \mathcal{P} and \mathbf{v} were arbitrary, this shows that $X^{\text{rd}}, \Xi \vdash \varphi^{\text{rd}}$.

Now suppose $X^{\text{rd}}, \Xi \vdash \varphi^{\text{rd}}$. Let $\mathcal{Q} = (\Omega^{\mathcal{Q}}, \Gamma^{\mathcal{Q}}, \mathbb{Q}^{\mathcal{Q}})$ be an \mathcal{L}' -model and let \mathbf{w} be an assignment into \mathcal{Q} such that $\mathcal{Q} \models \psi[\mathbf{w}]$ for all $\psi \in X \cup \Theta$. Since $\mathcal{Q} \models \Theta$, we may define \mathcal{P} as in Lemma 6.1.2. We then have $\mathcal{P} \models \Xi$ and $\mathcal{P}' \models \psi[\mathbf{v}']$ for all $\psi \in X \cup \Theta$. Proposition 6.1.3 then implies $\mathcal{P} \models \psi^{\text{rd}}[\mathbf{v}]$ for all $\psi \in X \cup \Theta$. In particular, $\mathcal{P} \models \psi[\mathbf{v}]$ for all $\psi \in X^{\text{rd}}$. Together with $\mathcal{P} \models \Xi$ and $X^{\text{rd}}, \Xi \vdash \varphi^{\text{rd}}$, this gives $\mathcal{P} \models \varphi^{\text{rd}}[\mathbf{v}]$. Another application of Proposition 6.1.3 gives $\mathcal{P}' \models \varphi[\mathbf{v}']$. Therefore, by Lemma 6.1.2, we have $\mathcal{Q} \models \varphi[\mathbf{w}]$. Since \mathcal{Q} and \mathbf{w} were arbitrary, this shows that $X, \Theta \vdash \varphi$. \square

$\langle \text{C:ded-elim-thm1} \rangle$ **Corollary 6.1.5.** *For any $X_1, X_2 \subseteq \mathcal{L}'$, we have $X_1 \cup \Theta \equiv X_2 \cup \Theta$ if and only if $X_1^{\text{rd}} \cup \Xi \equiv X_2^{\text{rd}} \cup \Xi$.*

Proof. Suppose $X_1 \cup \Theta \equiv X_2 \cup \Theta$. Then $X_1, \Theta \vdash X_2$. Theorem 6.1.4 implies $X_1^{\text{rd}}, \Xi \vdash X_2^{\text{rd}}$. Likewise, $X_2^{\text{rd}}, \Xi \vdash X_1^{\text{rd}}$. Therefore, $X_1^{\text{rd}} \cup \Xi \equiv X_2^{\text{rd}} \cup \Xi$. The proof of the converse is similar. \square

$\langle \text{D:ded-def-ext} \rangle$ **Definition 6.1.6.** Let $T \subseteq \mathcal{L}^0$ be a deductive theory, and let Θ and Ξ be as in Section 6.1.2. We say that Θ is *legitimate in T* if $\Xi \subseteq T$. If Θ is legitimate in T , then we define the deductive theory $T' \subseteq (\mathcal{L}')^0$ by $T' = T + \Theta$. The deductive theory T' is called a *definitorial extension of T* .

$\langle \text{C:ded-elim-thm2} \rangle$ **Corollary 6.1.7.** *Let T' be a definitorial extension of T . Then for all $\varphi \in \mathcal{L}'$, we have $T' \vdash \varphi$ if and only if $T \vdash \varphi^{\text{rd}}$. In particular, T is consistent if and only if T' is consistent.*

Proof. By Theorem 6.1.4, we have $T' \vdash \varphi$ if and only if $T^{\text{rd}}, \Xi \vdash \varphi^{\text{rd}}$. But $T^{\text{rd}} = T$ and $\Xi \subseteq T$. Therefore, $T^{\text{rd}}, \Xi \vdash \varphi^{\text{rd}}$ if and only if $T \vdash \varphi^{\text{rd}}$. \square

6.1.5 Inductive elimination

Our aim here is to prove Theorem 6.1.10 below, which is an inductive version of Theorem 6.1.4.

Let P be an inductive theory in \mathcal{L}^{IS} with root T_0 , and let T'_0 be a definitorial extension of T_0 . That is, $T'_0 = T_0 + \Theta$, where Θ is legitimate in T_0 , meaning that $\Xi \subseteq T_0$.

Let $X \subseteq (\mathcal{L}')^0$ and assume that $X \hookrightarrow T'_0$. Choose $\psi \in (\mathcal{L}')^0$ such that $X \equiv T'_0 + \psi$. Define $Q^X \subseteq (\mathcal{L}')^{\text{IS}}$ by

$$Q^X = \{(X, \varphi, p) \mid P(\varphi^{\text{rd}} \mid T_0, \psi^{\text{rd}}) = p\}.$$

We claim that the definition of Q^X does not depend on the choice ψ . To see this, suppose that $\zeta \in (\mathcal{L}')^0$ and $X \equiv T'_0 + \zeta$. Then $T'_0 \vdash \psi \rightarrow \zeta$. By Corollary 6.1.7, we have $T_0 \vdash \psi^{\text{rd}} \rightarrow \zeta^{\text{rd}}$, so that $T_0, \psi^{\text{rd}} \vdash \zeta^{\text{rd}}$. Reversing the roles of ψ and ζ gives $T_0, \zeta^{\text{rd}} \vdash \psi^{\text{rd}}$. Hence, $T_0 + \psi^{\text{rd}} = T_0 + \zeta^{\text{rd}}$. By the rule of logical equivalence, $P(\varphi^{\text{rd}} \mid T_0, \psi^{\text{rd}}) = p$ if and only if $P(\varphi^{\text{rd}} \mid T_0, \zeta^{\text{rd}}) = p$, and so Q^X does not depend on ψ .

Let $Q = \bigcup \{Q^X \mid X \hookrightarrow T'_0\}$. Then Q is strongly connected with root T'_0 .

(L:ind-elim-thm1) **Lemma 6.1.8.** *With notation as above, Q is satisfiable.*

Proof. Let \mathcal{P} be an \mathcal{L} -model such that $\mathcal{P} \models P$. Then $\mathcal{P} \models T_P$ and $\Xi \subseteq T_0 \subseteq T_P$, so that $\mathcal{P} \models \Xi$. Define \mathcal{P}' as in Proposition 6.1.3. Since $\mathcal{P}' \models \theta_s$ whenever $\mathcal{P} \models \xi_s$, we have $\mathcal{P}' \models \Theta$. Also, $T_0^{\text{rd}} = T_0$, so Proposition 6.1.3 gives $\mathcal{P}' \models T_0$. Therefore, $\mathcal{P}' \models T_0 + \Theta = T'_0$.

Now suppose $(X, \varphi, p) \in Q$. Write $X \equiv T'_0 + \psi$, where $P(\varphi^{\text{rd}} \mid T_0, \psi^{\text{rd}}) = p$. Then $\overline{\mathbb{P}} \varphi_{\Omega}^{\text{rd}} \cap \psi_{\Omega}^{\text{rd}} / \overline{\mathbb{P}} \psi_{\Omega}^{\text{rd}} = p$. By Proposition 6.1.3, we have $\overline{\mathbb{Q}} \varphi_{\Omega'} \cap \psi_{\Omega'} / \overline{\mathbb{Q}} \psi_{\Omega'} = p$, so that $\mathcal{P}' \models (X, \varphi, p)$. Since (X, φ, p) was arbitrary, $\mathcal{P}' \models Q$. \square

It follows from Lemma 6.1.8 and Theorem 4.2.7 that Q is consistent. We may therefore define $P' = \mathbf{P}_Q$. Let $P'_0 = P' \upharpoonright_{T'_0}$.

(L:ind-elim-thm2) **Lemma 6.1.9.** *With notation as above, $P'_0 = Q$.*

Proof. Since $Q \subseteq P'$ and $X \hookrightarrow T'_0$ for every $X \in \text{ante} Q$, we have $Q \subseteq P'_0$. Conversely, suppose that $P'_0(\varphi \mid X) = p$. Write $X \equiv T'_0 + \psi$, so that $P'(\varphi \mid T'_0, \psi) = p$. Let \mathcal{P} and \mathcal{P}' be as in the proof of Lemma 6.1.8. Since $\mathcal{P}' \models Q$, it follows that $\mathcal{P}' \models P'$. Therefore, $\overline{\mathbb{Q}} \varphi_{\Omega'} \cap \psi_{\Omega'} / \overline{\mathbb{Q}} \psi_{\Omega'} = p$. Proposition 6.1.3 then gives $\overline{\mathbb{P}} \varphi_{\Omega}^{\text{rd}} \cap \psi_{\Omega}^{\text{rd}} / \overline{\mathbb{P}} \psi_{\Omega}^{\text{rd}} = p$, so that $\mathcal{P} \models (T_0 + \psi^{\text{rd}}, \varphi^{\text{rd}}, p)$. Since this is true for every \mathcal{L} -model \mathcal{P} such that $\mathcal{P} \models P$, and since $T_0 + \psi^{\text{rd}} \hookrightarrow [T_0, T_P]$, it follows from Definition 4.2.10 that $P \models (T_0 + \psi^{\text{rd}}, \varphi^{\text{rd}}, p)$. Remark 4.2.15 therefore implies $P(\varphi^{\text{rd}} \mid T_0, \psi^{\text{rd}}) = p$, so that $(X, \varphi, p) \in Q$. Hence, $Q = P'_0$. \square

(T:ind-elim-thm) **Theorem 6.1.10 (Inductive elimination theorem).** *Let P be an inductive theory in \mathcal{L}^{IS} with root T_0 , and let T'_0 be a definitorial extension of T_0 as in Definition 6.1.6. Then there exists a unique inductive theory $P' \subseteq (\mathcal{L}')^{\text{IS}}$ with root T'_0 such that*

$$P'(\varphi \mid X, \Theta) = P(\varphi^{\text{rd}} \mid X^{\text{rd}}, \Xi), \quad (6.1.1) \quad \boxed{\text{ind-elim-thm}}$$

where either both sides exist or both sides do not.

Proof. Let P' be defined as above. We first show that $P'(\varphi \mid X, \Theta) = p$ implies $P(\varphi^{\text{rd}} \mid X^{\text{rd}}, \Xi) = p$. Suppose $P'(\varphi \mid X, \Theta) = p$. Write $X \cup \Theta \equiv T' + \psi$, where $T' \in [T'_0, T_{P'}]$, $\psi \in (\mathcal{L}')^0$, and $P'(\varphi \mid T'_0, \psi) = p$. Since $P'_0 = Q$, we have $Q(\varphi \mid T'_0, \psi) = p$. From the definition of Q , it follows that $P(\varphi^{\text{rd}} \mid T_0, \psi^{\text{rd}}) = p$. Therefore, to show that $P(\varphi^{\text{rd}} \mid X^{\text{rd}}, \Xi) = p$, it suffices to show that $X^{\text{rd}} \cup \Xi \equiv T + \psi^{\text{rd}}$ for some $T \in [T_0, T_P]$. For this, let $T = T((T')^{\text{rd}} \cup \Xi)$. Since $\Theta \subseteq T'_0 \subseteq T'$, we have $X \cup \Theta \equiv T' + \Theta + \psi$. Corollary 6.1.5 therefore gives $X^{\text{rd}} \cup \Xi \equiv (T')^{\text{rd}} \cup \Xi \cup \{\psi^{\text{rd}}\} \equiv T + \psi^{\text{rd}}$. We must now show that $T_0 \subseteq T$ and $T \subseteq T_P$.

For $T_0 \subseteq T$, note that $T_0 \subseteq T'_0 \subseteq T'$. We therefore have $T' \vdash T_0$, so that $T', \Theta \vdash T_0$. Since $T_0^{\text{rd}} = T_0$, Theorem 6.1.4 gives $(T')^{\text{rd}}, \Xi \vdash T_0$, so that $T_0 \subseteq T$.

For $T \subseteq T_P$, we first show that $T_{P'}^{\text{rd}} \subseteq T_P$. Let $\zeta \in T_{P'}$. Then $P'(\zeta \mid T'_0) = 1$. By Lemma 6.1.9, we have $Q(\zeta \mid T'_0) = 1$. From the definition of Q , it follows that $P(\zeta^{\text{rd}} \mid T_0) = 1$. Therefore, $\zeta^{\text{rd}} \in T_P$, and this shows that $T_{P'}^{\text{rd}} \subseteq T_P$.

Now, since $\Xi \subseteq T_0 \subseteq T_P$, we have that Θ is legitimate in T_P . Corollary 6.1.7 therefore gives $T_P + \Theta \vdash \zeta$ if and only if $T_P \vdash \zeta^{\text{rd}}$, for all $\zeta \in \mathcal{L}'$. By the above, $T_P \vdash T_{P'}^{\text{rd}}$. Hence, $T_P + \Theta \vdash T_{P'}$. Since $T' \subseteq T_{P'}$, this implies $T_P + \Theta \vdash T'$. Another application of Corollary 6.1.7 yields $T_P \vdash (T')^{\text{rd}}$. As previously noted, $\Xi \subseteq T_P$. Therefore, $T_P \vdash T$, so that $T \subseteq T_P$.

We now show that $P(\varphi^{\text{rd}} \mid X^{\text{rd}}, \Xi) = p$ implies $P'(\varphi \mid X, \Theta) = p$. Suppose $P(\varphi^{\text{rd}} \mid X^{\text{rd}}, \Xi) = p$. Write $X^{\text{rd}} \cup \Xi \equiv T + \psi$, where $T \in [T_0, T_P]$, $\psi \in \mathcal{L}^0$, and $P(\varphi^{\text{rd}} \mid T_0, \psi) = p$. By Lemma 6.1.9, the definition of Q , and the fact that $\psi^{\text{rd}} = \psi$, we have $P'(\varphi \mid T_0, \psi) = p$. Therefore, to show that $P'(\varphi \mid X, \Theta) = p$, it suffices to show that $X \cup \Theta \equiv T' + \psi$ for some $T' \in [T'_0, T_{P'}]$. For this, let $T' = T + \Theta$. Since $\Xi \subseteq T_0 \subseteq T$, we have $X^{\text{rd}} \cup \Xi \equiv T + \Xi + \psi = T' + \psi$.

We must now show that $T'_0 \subseteq T'$ and $T' \subseteq T_{P'}$. The first follows easily, since $T'_0 = T_0 + \Theta \subseteq T + \Theta = T'$. For the second, note that if $\zeta \in \mathcal{L}^0$ and $P(\zeta \mid T_0) = 1$, then $Q(\zeta \mid T'_0) = 1$, which implies $P'(\zeta \mid T'_0) = 1$. Therefore, $T_P \subseteq T_{P'}$. Since $T \subseteq T_P$, we have $T' = T + \Theta \subseteq T_P + \Theta \subseteq T_{P'} + \Theta$. But $\Theta \subseteq T'_0 \subseteq T_{P'}$, so it follows that $T' \subseteq T_{P'}$.

For uniqueness, let $P'' \subseteq (\mathcal{L}')^0$ be another inductive theory with root T'_0 that satisfies (6.1.1), and suppose $P''(\varphi \mid X) = p$. Then $\Theta \subseteq T'_0 \subseteq T(X)$, so that $X \equiv X \cup \Theta$. By the rule of logical equivalence and (6.1.1), we have

$$P''(\varphi \mid X) = P''(\varphi \mid X, \Theta) = P(\varphi^{\text{rd}} \mid X^{\text{rd}}, \Xi) = P'(\varphi \mid X, \Theta) = P'(\varphi \mid X),$$

which shows that $P'' \subseteq P'$. Reversing the roles of P' and P'' gives $P'' = P'$. \square

The inductive theory P' in Theorem 6.1.10 is called a *definitorial extension* of P . In the special case that P is a complete inductive theory, we have the following semantic characterization of P' .

(C:ind-elim-thm) **Corollary 6.1.11.** *Let $T_0 \subseteq \mathcal{L}^0$ be a deductive theory. Let \mathcal{P} be an \mathcal{L} -model with $\mathcal{P} \vDash T_0$ and define $P = \mathbf{Th} \mathcal{P} \downarrow_{[T_0, Th P]}$. Let P' be a definitorial extension of P . Then $P' = \mathbf{Th} \mathcal{P}' \downarrow_{[T'_0, Th \mathcal{P}]}$, where \mathcal{P}' is defined above Lemma 6.1.2.*

Proof. By Lemma 6.1.9 and Theorem 3.3.4, it suffices to show that $Q = \mathbf{Th} \mathcal{P}' \downarrow_{T'_0}$. Note that $(X, \varphi, p) \in Q$ if and only if we can write $X \equiv T'_0 + \psi$,

where $\mathcal{P} \models (T_0 + \psi^{\text{rd}}, \varphi^{\text{rd}}, p)$. Also note that $(X, \varphi, p) \in \mathbf{Th} \mathcal{P}' \downarrow_{T'_0}$ if and only if we can write $X \equiv T'_0 + \psi$, where $\mathcal{P}' \models (T'_0 + \psi, \varphi, p)$. Since P' is a definitorial extension, we have $\Xi \subseteq T_0$, which implies $\mathcal{P} \models \Xi$. Hence, from Proposition 6.1.3, it follows that $\mathcal{P} \models (T_0 + \psi^{\text{rd}}, \varphi^{\text{rd}}, p)$ if and only if $\mathcal{P}' \models (T'_0 + \psi, \varphi, p)$. \square

6.1.6 Primitive vs. defined symbols

(S:prim-vs-def)

Let $P \subseteq \mathcal{L}^{\text{IS}}$ be an inductive theory with root T_0 . The extralogical symbols in L have no formal definitions within P . Syntactically, they get their meaning from their use, that is, from the inductive statements $(X, \varphi, p) \in P$ in which they appear. If P and T_0 are generated by a set of axioms, then we might say that the symbols in L are “defined” by these axioms. That is, they get whatever meaning they may have from what the axioms have to say about them.

Now suppose P' is a definitorial extension of P . Then, unlike the symbols in L , the symbols in $L' \setminus L$ do have a formal definition in P' . For instance, to each constant symbol $c \in L' \setminus L$, there corresponds a defining formula $\delta_c(y) \in \mathcal{L}$ such that $\theta_c = \forall y(y = c \leftrightarrow \delta_c(y)) \in T'_0$. It seems, then, that the symbols in L' can be divided into two disjoint categories. Those in L we might call the “primitive symbols,” which lack a formal definition and get their meaning from their use. While those in $L' \setminus L$ we might call the “defined symbols,” which get their meaning by virtue of being defined in terms of the primitive symbols.

But this distinction is metatheoretical. Neither the deductive theory T'_0 nor the inductive theory P' can “see” which symbols are primitive and which are defined. From the point of view of P' , a constant symbol $c \in L' \setminus L$ is a primitive symbol that gets its meaning from the statements and sentences in P' and T'_0 , one of which is θ_c . In other words, there is no real difference between thinking of θ_c as a definition of c and thinking of θ_c as a new axiom that gives meaning to the primitive symbol c . In the context of a given inductive theory such as P' , there are no defined symbols. There are only primitive symbols.

Moreover, if we understand meaning as coming from use, which it does in formal logical systems such as this, then the very act of defining the new symbol $c \in L' \setminus L$ can change the meaning of the primitive symbols in L . More specifically, to define c , we must add the new “axiom,” $\theta_c = \forall y(y = c \leftrightarrow \delta_c(y))$, changing T_0 to T'_0 . The formula $\delta_c(y)$ undoubtedly uses symbols from L . Under the extension P' , those symbols are now used differently from how they were used under P . Hence, their meanings in P' and P may be different.

In practice, this latter issue may be one that we rarely, if ever, encounter. This is because we can avoid the issue in any circumstance where the elimination theorems apply. But we will later see in Chapter 7 (specifically Sections 7.3.4 and 7.5.3) circumstances where the elimination theorems do not apply. And in those circumstances, we must confront this issue head on.

6.2 Zermelo–Fraenkel set theory

(S:ZFC) In this section, we present Zermelo–Fraenkel set theory in the infinitary setting. Let \mathcal{L} be a language that contains a binary relation symbol \in . We use the boldface \in here to distinguish it from the usual \in that is used when discussing structures and models. As noted in Section 5.1.2, we will use the shorthand $(\forall y \in x)\varphi = \forall y(y \in x \rightarrow \varphi)$ and $(\exists y \in x)\varphi = \exists y(y \in x \wedge \varphi)$. More generally,

$$(\forall y_1 \cdots y_n \in x)\varphi = (\forall y_1 \in x) \cdots (\forall y_n \in x)\varphi,$$

and similarly for \exists . We also write $x \subseteq y = \forall z(z \in x \rightarrow z \in y)$. Again, we use boldface \subseteq to distinguish it from the usual \subseteq .

All of the formulas we define below are sentences. To simplify notation, we write them as open formulas of the form $\varphi(\vec{x})$. It is to be understood that this refers to the sentence $\forall \vec{x}\varphi(\vec{x})$.

6.2.1 Extensionality, union, and power set

Define the sentences

$$\begin{aligned} \text{AE} &: \forall z(z \in x \leftrightarrow z \in y) \rightarrow x = y \\ \text{AU} &: \forall x \exists y \forall z(z \in y \leftrightarrow (\exists u \in x)z \in u) \\ \text{AP} &: \forall x \exists y \forall z(z \in y \leftrightarrow z \subseteq x) \end{aligned}$$

These are, respectively, the axioms of extensionality, union, and power set. The first says that two sets x and y are equal if they have the same elements. The second says that, given a collection x of sets, there is a set y which is the union of the sets in x . The third says that if x is a set, then there is a set y consisting of all the subsets of x . We will have more to say about these axioms later.

6.2.2 Axiom schema of separation

(S:separation) For $\varphi(x, z, \vec{u}) \in \mathcal{L}$ with $y \notin \text{free } \varphi$, define

$$\text{AS}(\varphi) : \exists y \forall z(z \in y \leftrightarrow \varphi \wedge z \in x)$$

This is called the axiom of separation. Given a set x , the axiom $\text{AS}(\varphi)$ allows us to create the set $\{z \in x \mid \varphi(x, z, \vec{u})\}$, which depends on x and \vec{u} . This is done as follows. Let $\varphi(x, z, \vec{u}) \in \mathcal{L}$ with $y \notin \text{free } \varphi$. It can be shown that

$$\text{AE}, \text{AS}(\varphi) \vdash \forall x \vec{u} \exists! y \forall z(z \in y \leftrightarrow \varphi \wedge z \in x).$$

Hence, in any theory that contains AE and $\text{AS}(\varphi)$, we may make the legitimate definitorial extension, $y = Fx\vec{u} \leftrightarrow \forall z(z \in y \leftrightarrow \varphi \wedge z \in x)$. The term $Fx\vec{u}$ is exactly the set we were trying to create. In such a case, we use the notation $\{z \in x \mid \varphi\}$ or $\{z \in x \mid \varphi(x, z, \vec{u})\}$ as shorthand for the term $Fx\vec{u}$.

The axiom of separation is, in fact, an axiom schema. It is one axiom for each allowable formula φ . Let

$$\text{AS} = \{\text{AS}(\varphi) \mid \varphi(x, z, \vec{u}) \in \mathcal{L} \text{ and } y \notin \text{free } \varphi\},$$

and set $\text{AS}_{\text{fin}} = \text{AS} \cap \mathcal{L}_{\text{fin}}^0$. Note that $\text{AS}(\varphi) \in \mathcal{L}_{\text{fin}}^0$ if and only if $\varphi \in \mathcal{L}_{\text{fin}}$. Hence, AS_{fin} can be defined just as AS , but with the requirement that $\varphi \in \mathcal{L}_{\text{fin}}$. In other words, AS_{fin} is the usual axiom schema of separation used in first-order logic. The difference between AS and AS_{fin} is that with AS , we are allowed to use infinitary formulas φ when building new sets.

6.2.3 Axiom schema of replacement

(S:replacement) For $\varphi(x, y, \vec{z}) \in \mathcal{L}$ with $u, v \notin \text{free } \varphi$, define

$$\text{AR}(\varphi) : \forall x \exists! y \varphi \rightarrow \forall u \exists v \forall y (y \in v \leftrightarrow (\exists x \in u) \varphi)$$

This is called the axiom of replacement. Given a set u and a function F , it allows us to create a set of the form $\{Fx\vec{z} \mid x \in u\}$, which depends on u and \vec{z} . The function F is determined by the formula φ . This is done as follows. Let $\varphi(x, y, \vec{z}) \in \mathcal{L}$ with $u, v \notin \text{free } \varphi$. Suppose T is a theory such that $T \vdash \forall x \vec{z} \exists! y \varphi$. Also assume $\text{AE} \in T$ and $\text{AR}(\varphi) \in T$. In T , we may make the legitimate definitorial extension, $y = Fx\vec{z} \leftrightarrow \varphi(x, y, \vec{z})$. Let

$$\psi(u, v, \vec{z}) = \forall y (y \in v \leftrightarrow (\exists x \in u) \varphi),$$

and note that

$$\psi \equiv_T \forall y (y \in v \leftrightarrow (\exists x \in u) y = Fx\vec{z}).$$

By our hypotheses on T , it follows that $T \vdash \forall u \vec{z} \exists! v \psi(u, v, \vec{z})$. Hence, we may make the legitimate definitorial extension, $v = Gu\vec{z} \leftrightarrow \psi(u, v, \vec{z})$. The term $Gu\vec{z}$ is exactly the set we were trying to create. In such a case, we use the notation $\{Fx\vec{z} \mid x \in u\}$ as shorthand for the term $Gu\vec{z}$.

The axiom of replacement is, in fact, an axiom schema. It is one axiom for each allowable formula φ . Let

$$\text{AR} = \{\text{AR}(\varphi) \mid \varphi(x, y, \vec{z}) \in \mathcal{L} \text{ and } u, v \notin \text{free } \varphi\},$$

and set $\text{AR}_{\text{fin}} = \text{AR} \cap \mathcal{L}_{\text{fin}}^0$. Note that $\text{AR}(\varphi) \in \mathcal{L}_{\text{fin}}^0$ if and only if $\varphi \in \mathcal{L}_{\text{fin}}$. Hence, AR_{fin} can be defined just as AR , but with the requirement that $\varphi \in \mathcal{L}_{\text{fin}}$. In other words, AR_{fin} is the usual axiom schema of replacement used in first-order logic. The difference between AR and AR_{fin} is that with AR , we are allowed to use infinitary formulas to construct our defined function symbol F .

6.2.4 Definitorial extensions and shorthand

To state the remaining axioms, it will be useful to create new symbols, both shorthand and formally defined extralogical symbols. To this end, we define

$$\begin{aligned}\varphi_1(x, z) &: z \notin x \\ \varphi_2(x, y, z_1, z_2) &: (\forall u u \notin x) \wedge y = z_1 \vee (\exists u u \in x) \wedge y = z_2 \\ \varphi_3(x, z, u) &: z \in u\end{aligned}$$

Note that each φ_i is in \mathcal{L}_{fin} . Let T be the deductive theory generated by

$$\{\text{AE}, \text{AU}, \text{AP}, \text{AS}(\varphi_1), \text{AS}(\varphi_3), \text{AR}(\varphi_2)\}.$$

Since $\text{AE}, \text{AS}(\varphi_1) \vdash \exists! y \forall z z \notin y$, we may make the legitimate definitorial extension, $y = \emptyset \leftrightarrow \forall z z \notin y$. Letting

$$\varphi(x, y) = \forall z (z \in y \leftrightarrow (\exists u \in x) z \in u),$$

we have $\text{AE}, \text{AU} \vdash \forall x \exists! y \varphi(x, y)$, which allows the extension $y = \bigcup x \leftrightarrow \varphi(x, y)$. And with $\varphi(x, y) = \forall z (z \in y \leftrightarrow z \subseteq x)$, we have $\text{AE}, \text{AP} \vdash \forall x \exists! y \varphi(x, y)$, which allows the extension $y = \mathfrak{P}x \leftrightarrow \varphi(x, y)$.

We now have

$$\varphi_2 \equiv_T (x = \emptyset \wedge y = z_1 \vee x \neq \emptyset \wedge y = z_2).$$

It can be shown that $\text{AE}, \text{AS}(\varphi_1) \vdash \forall x \vec{z} \exists! y \varphi_2(x, y, \vec{z})$. This allows the extension $y = Fx\vec{z} \leftrightarrow \varphi_2(x, y, \vec{z})$. Hence, since $\text{AR}(\varphi_2) \in T$, we may adopt the shorthand $\{z_1, z_2\}$ for the term $\{Fx\vec{z} \mid x \in \mathfrak{P}\mathfrak{P}\emptyset\}$. We then write $\{z\}$ as shorthand for $\{z, z\}$. We can recursively define the shorthand $\{z_1, \dots, z_{n+1}\} = \{z_1, \dots, z_n\} \cup \{z_{n+1}\}$.

Since $\text{AE}, \text{AS}(\varphi_3) \in T$, we have the term $\{z \in x \mid \varphi_3\} = \{z \in x \mid z \in u\}$. We use the shorthand $x \cap u$ to denote this term. We also write $x \cup u$ as shorthand for $\bigcup\{x, u\}$, and $\mathbf{S}x$ as shorthand for the term $x \cup \{x\}$. Finally, the ordered pair (x, y) is defined via shorthand as the term $\{\{x\}, \{x, y\}\}$. We can recursively define the shorthand $(x_1, \dots, x_{n+1}) = ((x_1, \dots, x_n), x_{n+1})$.

6.2.5 Axioms of infinity, foundation, and choice

Now define the sentences

$$\begin{aligned}\text{AI} &: \exists u (\emptyset \in u \wedge \forall x (x \in u \rightarrow \mathbf{S}x \in u)) \\ \text{AF} &: (\forall x \neq \emptyset) (\exists y \in x) x \cap y = \emptyset \\ \text{AC} &: \forall u (\emptyset \notin u \wedge (\forall xy \in u) (x \neq y \rightarrow x \cap y = \emptyset) \rightarrow \exists z (\forall x \in u) \exists! y (y \in x \cap z))\end{aligned}$$

These are, respectively, the axioms of infinity, foundation, and choice. The axiom of infinity ensures the existence of an infinite set. The axiom of foundation, among other things, ensure that no set can be an element of itself. The axiom of choice asserts the existence of a set containing exactly one element from each of a given collection of sets.

6.2.6 Finitary and infinitary ZFC

We now define

$$\text{ZFC}_- = T + \text{AS}_{\text{fin}} + \text{AR}_{\text{fin}} + \{\text{AI}, \text{AF}, \text{AC}\},$$

and set $\text{ZFC}_{\text{fin}} = \text{ZFC}_- \cap \mathcal{L}_{\text{fin}}^0$. Note that $\text{ZFC}_- = T(\Lambda_-^{\text{ZFC}})$, where

$$\Lambda_-^{\text{ZFC}} = \{\text{AE}, \text{AU}, \text{AP}, \text{AI}, \text{AF}, \text{AC}\} \cup \text{AS}_{\text{fin}} \cup \text{AR}_{\text{fin}}$$

are the usual finitary axioms of set theory. Hence, from Proposition 5.3.15, it follows that ZFC_{fin} is the usual Zermelo–Fraenkel set theory with the axiom of choice, formulated in first-order logic.

We also define

$$\text{ZFC} = T + \text{AS} + \text{AR}_{\text{fin}} + \{\text{AI}, \text{AF}, \text{AC}\},$$

Note that $\text{ZFC} = T(\Lambda^{\text{ZFC}})$, where

$$\Lambda^{\text{ZFC}} = \{\text{AE}, \text{AU}, \text{AP}, \text{AI}, \text{AF}\} \cup \text{AS} \cup \text{AR}_{\text{fin}}.$$

The difference between ZFC and ZFC_- is that, in ZFC , we are allowed to use infinitary formulas when applying the axiom schema of separation.

Finally, we define

$$\text{ZFC}_+ = T + \text{AS} + \text{AR} + \{\text{AI}, \text{AF}, \text{AC}\},$$

Note that $\text{ZFC}_+ = T(\Lambda^{\text{ZFC}} \cup \text{AR})$. In ZFC_+ , we are also allowed to use infinitary formulas when applying the axiom schema of replacement. Also note that $\Lambda_-^{\text{ZFC}} \subseteq \Lambda^{\text{ZFC}}$ and $\text{ZFC}_{\text{fin}} \subseteq \text{ZFC}_- \subseteq \text{ZFC} \subseteq \text{ZFC}_+$.

By the same reasoning as in the proof of Proposition 5.3.22, we obtain the following.

(P:models-of-ZFC-) **Proposition 6.2.1.** *Let $\mathcal{P} = (\Omega, \Sigma, \mathbb{P})$ be a model. Then $\mathcal{P} \models \text{ZFC}_-$ if and only if $\omega \models \text{ZFC}_{\text{fin}}$ for \mathbb{P} -a.e. $\omega \in \Omega$. Consequently, ZFC_- is consistent if and only if ZFC_{fin} is consistent. Moreover, $\text{ZFC}_- \vdash \varphi$ if and only if $\omega \models \Lambda_-^{\text{ZFC}}$ implies $\omega \models \varphi[v]$ for all ω and all assignments v into ω .*

6.2.7 Consistency of ZFC

(S:Con-ZFC) In first-order logic, the consistency of ZFC_{fin} is implied by a variety of different sufficient conditions. One such condition is the following consequence of [6, Theorem 8.2.8].

(T:Con-ZFC-fin) **Theorem 6.2.2.** *If there exists a strongly inaccessible cardinal, then ZFC_{fin} is consistent in first-order logic.*

The existence of a strongly inaccessible cardinal cannot be proven in first-order logic using the standard axioms of set theory, Λ_-^{ZFC} . If it could, then these axioms could prove their own consistency, in violation of Gödel's second

incompleteness theorem (see [28, Theorem 7.3.2]). In situations where we want to use a strongly inaccessible cardinal, we must assume it exists, effectively adding its existence as a new axiom.

In the comprehensive articles, [23] and [24], Penelope Maddy discusses the metamathematical arguments, both historical and present, for accepting not only the current axioms, Λ^{ZFC} , but also additional possible axioms, including the assumption that strongly inaccessible cardinals exist. Of all the non-ZFC axioms discussed, this so-called *Axiom of Inaccessibles* seems to have the most support, with many historical backers, including Gödel. It is supported by a number of metamathematical principles, which she calls *maximize*, *inexhaustibility*, *uniformity*, *whimsical identity*, and *reflection*. This last principle, she says, is “probably the most universally accepted rule of thumb in higher set theory.”

In any case, if we want to have an inductive theory whose root contains either ZFC_- or ZFC , then we will need to assume something that ensures these theories are consistent. Of all the assumptions we could make in this regard, we prefer the one in Theorem 6.2.2. It is a very mild assumption, compared to others we might make. It has good metamathematical support. And it seems to produce a number of helpful results for us. Our first example of this is Theorem 6.2.4 below, which shows that this same assumption gives us the consistency of ZFC .

To prove this result, we must first see how the existence of a strongly inaccessible cardinal implies the consistency of ZFC_{fin} in Theorem 6.2.2. For this, we begin by defining the *von Neumann hierarchy*, which is a collection of sets indexed by the ordinals. For each ordinal α , we recursively define the set V_α as follows. Let $V_0 = \emptyset$. If $\alpha = \beta + 1$, then let $V_{\beta+1} = \mathfrak{P}V_\beta$. If α is a limit ordinal, then let $V_\alpha = \bigcup_{\xi < \alpha} V_\xi$. The sets V_α satisfy $V_\alpha = \bigcup_{\beta < \alpha} \mathfrak{P}V_\beta$, for any ordinal α . Each V_α is a transitive set, meaning that if $A \in V_\alpha$ and $x \in A$, then $x \in V_\alpha$. The sets V_α also satisfy the following properties:

- (i) $\beta < \alpha$ implies $V_\beta \in V_\alpha$ and $V_\beta \subseteq V_\alpha$,
- (ii) $A \in V_\alpha$ implies $A \subseteq V_\alpha$,
- (iii) $B \subseteq V_\beta$ implies $B \in V_\alpha$ for all $\alpha > \beta$, and
- (iv) $\alpha \subseteq V_\alpha$ for all α .

Using our identification of natural numbers and ordinals, the first five sets in the von Neumann hierarchy can be written as

$$\begin{aligned}
 V_0 &= \emptyset, \\
 V_1 &= \{0\}, \\
 V_2 &= \{0, 1\}, \\
 V_3 &= \{0, 1, 2, \{1\}\}, \\
 V_4 &= \{0, 1, 2, 3, \{1\}, \{2\}, \{0, 2\}, \{1, 2\}, \{\{1\}\}, \{0, \{1\}\}, \{0, 1, \{1\}\}, \{0, 1, 2, \{1\}\}, \\
 &\quad \{1, \{1\}\}, \{2, \{1\}\}, \{0, 2, \{1\}\}, \{1, 2, \{1\}\}\}.
 \end{aligned}$$

Note that $|V_5| = 65536$ and $|V_6| = 2^{65536}$. Although these sets grow rapidly, we have $|V_n| < \infty$ for all $n \in \mathbb{N}_0$. Also note that $n \subseteq V_n$ and $n \in V_{n+1}$, for all $n \in \mathbb{N}_0$. Since $\mathbb{N}_0 = \omega$, we have $\mathbb{N}_0 \subseteq V_\omega$ and $\mathbb{N}_0 \in V_{\omega+1}$.

If κ is a cardinal number, then we define the L -structure $\nu_\kappa = (V_\kappa, \in^{\nu_\kappa})$ by setting $\in^{\nu_\kappa} = \in$.

$\langle \text{L:V-is-closed} \rangle$ **Lemma 6.2.3.** *If $\nu_\kappa \models \Lambda_-^{\text{ZFC}}$, then V_κ satisfies the following:*

- (i) $\nu_\kappa \models (\forall x x \notin y)[b]$ if and only if $b = \emptyset$,
- (ii) $\nu_\kappa \models (x \subseteq y)[a, b]$ if and only if $a \subseteq b$,
- (iii) if $a \subseteq b$ and $b \in V_\kappa$, then $a \in V_\kappa$, and
- (iv) if $b \in V_\kappa$, then $\mathfrak{P}b \in V_\kappa$.

Proof. Let $b \in V_\kappa$. Note that $\nu_\kappa \models (\forall x x \notin y)[b]$ if and only if, for all $a \in V_\kappa$, we have $a \notin b$. Suppose $b \neq \emptyset$. Then there is an object a such that $a \in b$. Since V_κ is a transitive set, this implies $a \in V_\kappa$, a contradiction. Hence, (i) holds.

Next, let $a, b \in V_\kappa$. Note that $\nu_\kappa \models (x \subseteq y)[a, b]$ if and only if, for all $c \in V_\kappa$, we have $c \in a$ implies $c \in b$. Suppose $a \not\subseteq b$. Then there is an object c such that $c \in a$ but $c \notin b$. Again, since V_κ is transitive, we have $c \in V_\kappa$, a contradiction. Therefore, (ii) holds.

Let $a \subseteq b$ and $b \in V_\kappa$. Since every infinite cardinal number is a limit ordinal, we have $V_\kappa = \bigcup_{\beta < \kappa} V_\beta$ and $V_\kappa = \bigcup_{\beta < \kappa} \mathfrak{P}V_\beta$. Choose $\beta < \kappa$ such that $b \in V_\beta$, and note that $\beta + 1 < \kappa$. Then $b \subseteq V_\beta$, so that $a \subseteq V_\beta$, which implies $a \in V_{\beta+1} \subseteq V_\kappa$, and (iii) holds.

Finally, suppose $b \in V_\kappa$. Choose $\beta < \kappa$ such that $b \in V_\beta$. Then $b \subseteq V_\beta$, so that every $a \subseteq b$ satisfies $a \in V_{\beta+1}$. Hence, $\mathfrak{P}b \subseteq V_{\beta+1}$, and this implies that $\mathfrak{P}b \in V_{\beta+2} \subseteq V_\kappa$. \square

$\langle \text{T:Con-ZFC} \rangle$ **Theorem 6.2.4.** *Let ν_κ be as above. Then the following are equivalent:*

- (i) κ is strongly inaccessible,
- (ii) $\nu_\kappa \models \Lambda_-^{\text{ZFC}}$,
- (iii) $\nu_\kappa \models \Lambda^{\text{ZFC}}$

In particular, if there exists a strongly inaccessible cardinal, then ZFC is strictly satisfiable.

Proof. The equivalence of (i) and (ii) is [6, Theorem 8.2.8], and (iii) implies (ii) follows from the fact that $\Lambda_-^{\text{ZFC}} \subseteq \Lambda^{\text{ZFC}}$. For the final implication, assume (ii) holds. We need to show that, for all $\varphi \in \mathcal{L}$,

$$\text{if } \varphi = \varphi(x, z, \vec{u}) \text{ and } y \notin \text{free } \varphi, \text{ then } \nu_\kappa \models \text{AS}(\varphi). \quad (6.2.1) \quad \boxed{\text{Con-ZFC}}$$

We will prove this by induction on φ . Since prime formulas are finitary and $\nu_\kappa \models \Lambda_-^{\text{ZFC}}$, it holds whenever φ is prime.

Suppose $\varphi = \neg\psi$ and (6.2.1) holds for ψ . Assume $\varphi = \varphi(x, z, \vec{u})$ and $y \notin \text{free } \varphi$. Then the same is true for ψ . Hence, $\nu_\kappa \models \text{AS}(\psi)$. Let $a, d_1, \dots, d_n \in V_\kappa$. Then there exists $b \in V_\kappa$ such that, for all $c \in V_\kappa$, we have $c \in b$ if and only if $c \in a$ and $\nu_\kappa \models \psi[a, c, \vec{d}]$. Since V_κ is transitive, this implies $b = \{c \in a \mid \nu_\kappa \models \psi[a, c, \vec{d}]\}$. By Lemma 6.2.3(iii), we have $a \setminus b \in V_\kappa$. But $a \setminus b = \{c \in a \mid \nu_\kappa \models (\neg\psi)[a, c, \vec{d}]\}$. Hence, $\nu_\kappa \models \text{AS}(\neg\psi)$.

Now suppose $\varphi = \bigwedge \Phi$ and (6.2.1) holds for each $\theta \in \Phi$, and assume that $\varphi = \varphi(x, z, \vec{u})$ and $y \notin \text{free } \varphi$. Then the same is true for each θ , so that $\nu_\kappa \models \text{AS}(\theta)$. Let $a, d_1, \dots, d_n \in V_\kappa$. As above, $b_\theta = \{c \in a \mid \nu_\kappa \models \theta[a, c, \vec{d}]\} \in V_\kappa$. Lemma 6.2.3(iii) then gives $b = \bigcap_{\theta \in \Phi} b_\theta \in V_\kappa$. But $b = \{c \in a \mid \nu_\kappa \models (\bigwedge \Phi)[a, c, \vec{d}]\}$, so that $\nu_\kappa \models \text{AS}(\bigwedge \Phi)$.

Finally, suppose $\varphi = \forall v\psi$ and (6.2.1) holds for ψ . Assume $\varphi = \varphi(x, z, \vec{u})$ and $y \notin \text{free } \varphi$. Rename v to u_{n+1} , so that $\nu_\kappa \models \text{AS}(\psi)$. Let $a, d_1, \dots, d_n \in V_\kappa$. As above, for each $e \in V_\kappa$, we have $b_e = \{c \in a \mid \nu_\kappa \models \psi[a, c, \vec{d}, e]\} \in V_\kappa$. By Lemma 6.2.3(iii), it follows that $b = \bigcap_{e \in V_\kappa} b_e \in V_\kappa$. But $b = \{c \in a \mid \nu_\kappa \models (\forall v\psi)[a, c, \vec{d}]\}$, so that $\nu_\kappa \models \text{AS}(\forall v\psi)$. \square

6.3 Real inductive theories in ZFC₋

(S:ind-th-ZFC₋) In this section, we show how to represent real numbers in ZFC₋. We then use this representation to construct real inductive theories in ZFC₋.

6.3.1 The set of natural numbers

Let us add to the language of ZFC₋ a constant symbol \underline{n} for each $n \in \mathbb{N}_0$. We do this through the definitorial extensions, $y = \underline{0} \leftrightarrow y = \emptyset$ and $y = \underline{n} \leftrightarrow y = \mathbf{S} \cdots \mathbf{S}\emptyset$, where \mathbf{S} is repeated n times. The symbols \underline{n} are syntactic representations of the natural numbers. They give us a way to define each individual natural number in ZFC₋. But defining each individual natural number is not the same as defining the set of natural numbers. We must define a set that contains the natural numbers, and nothing but the natural numbers. One particularly counterintuitive result in ZFC₋ is that this is impossible, in the sense expressed by (6.3.1) below.

Instead, the best we can do in ZFC₋ is to prove that there is a smallest set that contains the natural numbers. But we cannot prove that this smallest set contains only the natural numbers. For instance, it is consistent with ZFC₋ to postulate that this smallest set contains an object which is greater than every natural number. Such an object would be called a *nonstandard natural number*. Similarly, ZFC₋ is consistent with the existence of nonstandard integers, rationals, and real numbers. On the other hand, in ZFC, we can prove that there is a unique set that contains the natural numbers, and nothing else. This is expressed in (6.3.2) below, whose proof uses essentially the same technique as in the proof of Proposition 5.3.23.

In ZFC₋, we define the smallest set that contains the natural numbers as

follows. Let $\varphi(u) = \underline{0} \in u \wedge \forall x(x \in u \rightarrow \mathbf{S}x \in u)$, so that $\mathbf{AI} = \exists u \varphi(u)$, and let $\delta(y) = \varphi(y) \wedge \forall z(\varphi(z) \rightarrow y \subseteq z)$.

(P: defn-N-ZFC) **Proposition 6.3.1.** *With notation as above, we have $\Lambda_-^{\text{ZFC}} \vdash \exists! y \delta(y)$.*

Proof. Let $\psi(x) = \forall z(\varphi(z) \rightarrow x \in z)$. Since $ZFC_- \vdash \mathbf{AS}(\psi)$, we may define the term $t(u) = \{x \in u \mid \psi(x)\}$, which satisfies $ZFC_- \vdash \varphi(u) \rightarrow \delta(t(u))$. Hence, since $ZFC_- \vdash \mathbf{AI}$ and $\mathbf{AI} = \exists u \varphi(u)$, it follows that $ZFC_- \vdash \exists y \delta(y)$. For uniqueness, simply note that $ZFC_- \vdash \delta(x) \wedge \delta(y) \rightarrow x \subseteq y \wedge y \subseteq x$. \square

By Proposition 6.3.1, we can make the legitimate definitorial extension $y = \underline{\mathbb{N}}_0 \leftrightarrow \delta(y)$. In words, $\underline{\mathbb{N}}_0$ is the smallest set that contains $\underline{0}$ and is closed under the successor operation.

6.3.2 Arithmetic operations

Recall that the ordered pair (x, y) is shorthand for the set $\{\{x\}, \{x, y\}\}$. Note that if $ZFC_- \vdash x \in u, y \in v$, then $ZFC_- \vdash (x, y) \in \mathfrak{P} \mathfrak{P}(u \cup v)$. Hence, we may adopt the shorthand

$$u \times v = \{z \in \mathfrak{P} \mathfrak{P}(u \cup v) \mid \exists xy(z = (x, y) \wedge x \in u \wedge y \in v)\},$$

and

$$v^u = \{z \in \mathfrak{P}(u \times v) \mid (\forall x \in u)(\exists! y \in v)(x, y) \in z\}.$$

Here, $u \times v$ is the Cartesian product of u and v , and v^u is the set of functions from u to v . We adopt the usual shorthand function notation, $z(x)$. Namely, if $\varphi(y) \in \mathcal{L}$, then $\varphi(z(x))$ is shorthand for the sentence,

$$\exists w(z \in v^u \wedge x \in u) \wedge \exists y((x, y) \in z \wedge \varphi(y)).$$

In other words, the sentence $\varphi(z(x))$ says that z is a function, x is in the domain of z , the object y is the unique object such that (x, y) is in z , and $\varphi(y)$ holds.

With these notions in hand, we can define addition in one of several equivalent ways, each ending with an explicitly defined constant symbol $+$ such that $ZFC_- \vdash + \in \underline{\mathbb{N}}_0^{\underline{\mathbb{N}}_0 \times \underline{\mathbb{N}}_0}$, and which agrees with ordinary addition on \mathbb{N}_0 . This latter fact means, specifically, that

$$ZFC_- \vdash (\underline{m}, \underline{n}, \underline{k}) \in + \text{ if and only if } m + n = k.$$

Note that in the above, the first instance of $+$ is an extralogical symbol and the second denotes ordinary addition of natural numbers. Since we use the same typographical symbol for both, we will need to rely on context to tell the difference. We adopt the usual shorthand, writing $x + y = z$ to mean $(x, y, z) \in +$. We then do the same for \cdot and $<$.

6.3.3 Peano arithmetic and nonstandard numbers

The axioms of Peano arithmetic can all be translated into the language of ZFC_- by replacing each quantifier $\forall x$ with $(\forall x \in \underline{\mathbb{N}}_0)$. By an abuse of notation, we will denote these translated axioms also by Λ_-^{PA} . We perform a similar abuse with Λ^{PA} , PA_- , and PA . It can be shown that $\text{PA}_- \subseteq \text{ZFC}_-$, and in a similar way, that $\text{PA} \subseteq \text{ZFC}$. Combined with Proposition 5.3.23, this latter fact tells us that if $\mathcal{N} \models \varphi$, then $\text{ZFC} \vdash \varphi$.

The arithmetical completeness of ZFC is intrinsically connected to the following result. The first part of the result, (6.3.1), shows the impossibility of proving in ZFC_- that $\underline{0}, \underline{1}, \underline{2}, \dots$ are the only natural numbers. The second part, (6.3.2), shows that this is not a problem in ZFC .

Proposition 6.3.2. *If ZFC_- is consistent, then*

$$\text{ZFC}_- \not\vdash \exists y \forall x (x \in y \leftrightarrow \bigvee_{n \in \mathbb{N}_0} x = \underline{n}). \quad (6.3.1) \text{nonstd-N-1}$$

On the other hand,

$$\text{ZFC} \vdash \forall x (x \in \underline{\mathbb{N}}_0 \leftrightarrow \bigvee_{n \in \mathbb{N}_0} x = \underline{n}). \quad (6.3.2) \text{nonstd-N-2}$$

Proof. For notational simplicity, let $\psi(y) = \forall x (x \in y \leftrightarrow \bigvee_{n \in \mathbb{N}_0} x = \underline{n})$, so that we aim to show $\text{ZFC}_- \not\vdash \exists y \psi(y)$ and $\text{ZFC} \vdash \psi(\underline{\mathbb{N}}_0)$.

Suppose ZFC_- is consistent and $\text{ZFC}_- \vdash \exists y \psi(y)$. Let $\varphi(u)$ and $\delta(y)$ be as in the definition of $\underline{\mathbb{N}}_0$. Note that $\text{ZFC}_- \vdash \psi(y) \rightarrow \varphi(y)$. Also, $\text{ZFC}_- \vdash \psi(y) \rightarrow \varphi(z) \rightarrow y \subseteq z$. Hence, $\text{ZFC}_- \vdash \psi(y) \rightarrow \delta(y)$. It therefore follows from $\text{ZFC}_- \vdash \exists y \psi(y)$ that $\text{ZFC}_- \vdash \psi(\underline{\mathbb{N}}_0)$. In particular, we have $\text{ZFC}_- \vdash \zeta(x)$, where $\zeta(x) = x \in \underline{\mathbb{N}}_0 \rightarrow \bigvee_{n \in \mathbb{N}_0} x = \underline{n}$.

By Proposition 6.2.1, since ZFC_- is consistent, we may find a structure ω such that $\omega \models \Lambda_-^{\text{ZFC}}$. Let c be a constant not in L and define $X \subseteq (\mathcal{L}c)_{\text{fin}}^0$ by $X = \Lambda_-^{\text{ZFC}} \cup \{c \in \underline{\mathbb{N}}_0\} \cup \{c \neq \underline{n} \mid n \in \mathbb{N}_0\}$. Let $X_0 \subseteq X$ be finite. Choose $m \in \mathbb{N}_0$ such that

$$X_0 \subseteq \Lambda_-^{\text{ZFC}} \cup \{c \in \underline{\mathbb{N}}_0\} \cup \{c \neq \underline{n} \mid n < m\}.$$

Let ω' be the Lc -expansion of ω with $c^{\omega'} = \underline{m}$. Then $\omega' \models X_0$, so that X_0 is strictly satisfiable. By the finiteness theorem for first-order logic (see, for example, [28, Theorem 3.3.1]), the set X is strictly satisfiable. Choose an Lc -structure ν such that $\nu \models X$, and let ν_0 be its L -reduct. Let $a = c^\nu$. Then $\nu_0 \models \Lambda_-^{\text{ZFC}}$, $\nu_0 \models (x \in \underline{\mathbb{N}}_0)[a]$, and $\nu_0 \models (x \neq \underline{n})[a]$ for all $n \in \mathbb{N}_0$. Therefore, $\nu_0 \not\models \zeta[a]$. Now let $\mathcal{P} = (\{\nu_0\}, \mathfrak{P}\{\nu_0\}, \delta_{\nu_0})$ and $\mathbf{v} = \langle v \rangle$, where $v(x) = a$. Then $\mathcal{P} \models \text{ZFC}_-$ and $\mathcal{P} \not\models \zeta[\mathbf{v}]$. By Theorem 5.3.21, we have $\text{ZFC}_- \not\vdash \zeta(x)$, a contradiction. This proves (6.3.1).

We now consider (6.3.2). By the definition of $\underline{\mathbb{N}}_0$, we have $\text{ZFC}_- \vdash \bigvee_{n \in \mathbb{N}_0} x = \underline{n} \rightarrow x \in \underline{\mathbb{N}}_0$. Therefore, we need only prove that $\text{ZFC} \vdash \zeta(x)$. By Definition 5.2.5(vii)', it suffices to show

$$\text{ZFC} \vdash \forall x \zeta = (\forall x \in \underline{\mathbb{N}}_0) \xi(x), \quad (6.3.3) \text{nonstd-N-3}$$

where $\xi(x) = (\bigvee_{n \in \mathbb{N}_0} x = \underline{n})$.

By $AS(\xi)$, we may define the term $t = \{x \in \underline{\mathbb{N}}_0 \mid \xi(x)\}$. It is straightforward to verify that $ZFC \vdash \varphi(t)$ and $ZFC \vdash \varphi(z) \rightarrow t \subseteq z$. Hence, $ZFC \vdash t = \underline{\mathbb{N}}_0$. Since $t = \underline{\mathbb{N}}_0 \equiv_{ZFC_-} (\forall x \in \underline{\mathbb{N}}_0) \xi(x)$, this gives (6.3.3). \square

6.3.4 Real numbers in ZFC_-

In our construction of ZFC_- in Section 6.2, we implemented a number of definitorial extensions, so that the extralogical signature of \mathcal{L} contains not only \in , but many other explicitly defined symbols, such as \emptyset , \cup , and \mathfrak{P} . For the present treatment, we consider ZFC_- to begin with the extralogical signature,

$$L = \{\in, \emptyset, \mathfrak{S}, \underline{\mathbb{N}}_0, +, \cdot, <\} \cup \{\underline{n} \mid n \in \mathbb{N}_0\}.$$

Any symbol which was defined in Section 6.2 but does not appear above is considered to be reduced. We may continue to use some of those symbols, but these uses should be considered shorthand, until otherwise specified. Note that each symbol in L , except for \in , is a constant symbol, and each is explicitly defined with a finitary defining formula.

From this starting point, we can now explicitly define each integer $z \in \mathbb{Z}$ in the usual way as an equivalence class of ordered pairs of natural numbers. For example, -5 is explicitly defined by the formula,

$$\delta(y) = \forall x(x \in y \leftrightarrow (\exists uv \in \underline{\mathbb{N}}_0)(x = (u, v) \wedge u + \underline{5} = v)).$$

The set of such equivalence classes is also explicitly definable, so that we may add the symbol $\underline{\mathbb{Z}}$. However, as was the case with $\underline{\mathbb{N}}_0$, we have that $ZFC_- \not\vdash \forall x(x \in \underline{\mathbb{Z}} \leftrightarrow \bigvee_{z \in \mathbb{Z}} x = \underline{z})$, provided ZFC_- is consistent.

We can then explicitly define $+_{\mathbb{Z}}$, $\cdot_{\mathbb{Z}}$, and $<_{\mathbb{Z}}$ for integers. In this way we obtain a definitorial extension of ZFC_- with signature

$$L = \{\in, \emptyset, \mathfrak{S}, \underline{\mathbb{N}}_0, \underline{\mathbb{Z}}, +, \cdot, <, +_{\mathbb{Z}}, \cdot_{\mathbb{Z}}, <_{\mathbb{Z}}\} \cup \{\underline{n} \mid n \in \mathbb{N}_0\} \cup \{\underline{z} \mid z \in \mathbb{Z}\},$$

where each symbol in L , except for \in , is a constant symbol, and each is explicitly defined with a finitary defining formula. We will omit the duplicates of $+$, \cdot , and $<$, and leave the distinction to context. Similarly, we will omit $\{\underline{n} \mid n \in \mathbb{N}_0\}$, and leave to context the distinction between the natural number n and the integer n .

Finally, we do the same for each $q \in \mathbb{Q}$ and for \mathbb{Q} itself, giving us the signature,

$$L = \{\in, \emptyset, \mathfrak{S}, \underline{\mathbb{N}}_0, \underline{\mathbb{Z}}, \underline{\mathbb{Q}}, +, \cdot, <\} \cup \{q \mid q \in \mathbb{Q}\}.$$

Again, everything but \in is a constant symbol with a finitary definition.

A set B is a *Dedekind cut* if B is a nonempty, proper subset of \mathbb{Q} that is downward closed and has no maximum element. Let

$$\begin{aligned} \varphi_{DC}(u) = & u \in \mathfrak{P}\underline{\mathbb{Q}} \wedge u \neq \emptyset \wedge u \neq \underline{\mathbb{Q}} \\ & \wedge (\forall x \in u)(y \in \underline{\mathbb{Q}} \wedge y < x \rightarrow y \in u) \wedge (\forall x \in u)(\exists y \in u)(x < y). \end{aligned}$$

Then $\varphi_{DC}(u)$ says that u is a Dedekind cut. Since φ_{DC} is finitary, we have $\text{ZFC}_- \vdash \text{AS}(\varphi_{DC})$. Hence, if

$$\delta(y) = \forall x(x \in y \leftrightarrow x \in \mathfrak{P}\underline{\mathbb{Q}} \wedge \varphi_{DC}(x)),$$

then $\text{ZFC}_- \vdash \exists!y \delta(y)$. We may therefore explicitly define $\underline{\mathbb{R}}$ by $y = \underline{\mathbb{R}} \leftrightarrow \delta(y)$. We can also explicitly define addition, multiplication, and less-than in $\underline{\mathbb{R}}$, all with finitary formulas. This gives us the extralogical signature,

$$L_- = \{\in, \emptyset, \mathbf{S}, \underline{\mathbb{N}}_0, \underline{\mathbb{Z}}, \underline{\mathbb{Q}}, \underline{\mathbb{R}}, +, \cdot, <\} \cup \{q \mid q \in \underline{\mathbb{Q}}\}.$$

As before, we have omitted the duplicate versions of $+$, \cdot , and $<$, and will rely on context to understand them. Also, each symbol in L_- , except for \in , is an explicitly defined constant symbol with a finitary defining formula.

6.3.5 The standard real structure

$\langle \text{S:std-R} \rangle$ Before constructing real inductive theories in ZFC_- , we first show a simpler, albeit more limited approach. This approach is essentially just a special case of Theorem 5.4.2.

Let $L_{\mathbb{R}} = \{+, \cdot, <\} \cup \{r \mid r \in \mathbb{R}\}$. In $\mathcal{L}_{\mathbb{R}}$, we will write $x \leq y$ as shorthand for $x < y \vee x = y$. Define the $L_{\mathbb{R}}$ -structure $\mathcal{R} = (\mathbb{R}, L^{\mathcal{R}})$ by letting $+^{\mathcal{R}}$, $\cdot^{\mathcal{R}}$, and $<^{\mathcal{R}}$ denote their ordinary counterparts in \mathbb{R} , and setting $r^{\mathcal{R}} = r$. Define the deductive theory $T_{\mathbb{R}}$ by $T_{\mathbb{R}} = \{\varphi \in \mathcal{L}^0 \mid \mathcal{R} \models \varphi\}$.

Let L be an extralogical signature with $L_{\mathbb{R}} \subseteq L$. If $P \subseteq \mathcal{L}^{\text{IS}}$ is an inductive theory with root $T_0 \supseteq T_{\mathbb{R}}$, then P is called a *real inductive theory* in $T_{\mathbb{R}}$.

In Theorem 5.4.2, we saw that every measure-theoretic probability model can be represented by an inductive model in a certain language. In Theorem 6.3.3 below, we show that if that measure-theoretic model is real-valued, then it can be represented by an inductive model in $\mathcal{L}_{\mathbb{R}}$.

Let (S, Γ, ν) be a probability space and let $X = \langle X_i \mid i \in I \rangle$ be an indexed collection of real-valued random variables. That is, each X_i is a Borel-measurable function from S to \mathbb{R} . Assume $\Gamma = \sigma(\langle X_i \mid i \in I \rangle)$.

Let $C = \{\underline{X}_i \mid i \in I\}$ be a set of distinct constant symbols not in $L_{\mathbb{R}}$, and define $L = L_{\mathbb{R}}C$.

$\langle \text{T:prob-model-iso-TR} \rangle$ **Theorem 6.3.3.** *There exists an \mathcal{L} -model $\mathcal{P} = (\Omega, \Sigma, \mathbb{P})$ with $\mathcal{P} \models T_{\mathbb{R}}$, and a function $h : S \rightarrow \Omega$ mapping $x \in S$ to $\omega \in \Omega$ such that*

- (i) $x \in \{X_i \leq r\}$ if and only if $\omega \models (\underline{X}_i \leq r)$,
- (ii) each $U \in \Gamma$ can be written as $U = h^{-1}\varphi_{\Omega}$ for some $\varphi \in \mathcal{L}^0$, and
- (iii) h induces a measure-space isomorphism from (S, Γ, ν) to \mathcal{P} .

Consequently, if $P = \mathbf{Th} \mathcal{P} \upharpoonright_{[T_{\mathbb{R}}, Th \mathcal{P}]}$, then

$$P(\bigwedge_{k=1}^n \underline{X}_{i(k)} \leq r_k \mid T_{\mathbb{R}}) = \nu \bigcap_{k=1}^n \{X_{i(k)} \leq r_k\}, \quad (6.3.4) \quad \boxed{\text{prob-model-iso-TR}}$$

whenever $i(1), \dots, i(n) \in I$ and $r_k \in \mathbb{R}$.

Proof. For each $x \in S$, define $\omega = \omega^x$ to be the L -expansion of \mathcal{R} given by $\omega^{\underline{X}_i} = X_i(x)$. Let $\Omega = \{\omega^x \mid x \in S\}$ and let $h : S \rightarrow \Omega$ denote the map $x \mapsto \omega^x$. Let $\mathcal{P} = (\Omega, \Sigma, \mathbb{P})$ be the measure space image of (S, Γ, ν) under h . Since $\omega \models T_{\mathbb{R}}$ for all $\omega \in \Omega$, we have $\mathcal{P} \models T_{\mathbb{R}}$. By construction, we have $X_i(x) \leq r$ if and only if $\omega^x \models \underline{X}_i \leq r$, so (i) holds.

For (ii), let

$$\Gamma' = \{U \in \Gamma \mid U = h^{-1}\varphi_{\Omega} \text{ for some } \varphi \in \mathcal{L}^0\} \subseteq \Gamma.$$

Since $\bigcup_n h^{-1}(\varphi_n)_{\Omega} = h^{-1}(\bigvee_n \varphi_n)_{\Omega}$ and $\perp_{\Omega} = \emptyset$, we have that Γ' is a σ -algebra. Let $r \in \mathbb{R}$. Since $X_i(x) \leq r$ if and only if $\omega^x \models \underline{X}_i \leq r$, it follows that $\{X_i \leq r\} = h^{-1}(\underline{X}_i \leq r)_{\Omega}$. Thus, $\{X_i \leq r\} \in \Gamma'$, so that X_i is Γ' -measurable for all $i \in I$. Since Γ is the smallest σ -algebra with this property, we have $\Gamma \subseteq \Gamma'$. Hence, $\Gamma = \Gamma'$, so (ii) holds.

If $U \in \Gamma$, $\varphi \in \mathcal{L}^0$, and $U = h^{-1}\varphi_{\Omega}$, then by the construction of \mathcal{P} , we have $\varphi_{\Omega} \in \Sigma$. Therefore, (ii) implies (iii).

Finally, since h also induces an isomorphism from $(S, \bar{\Gamma}, \bar{\nu})$ to $(\Omega, \bar{\Sigma}, \bar{\mathbb{P}})$, we have $\bar{\mathbb{P}} = \bar{\nu} \circ h^{-1}$. This gives

$$P(\bigwedge_{k=1}^n \underline{X}_{i(k)} \leq \underline{r}_k \mid T_{\mathbb{R}}) = \bar{\mathbb{P}} \bigcap_{k=1}^n (\underline{X}_{i(k)} \leq \underline{r}_k)_{\Omega} = \nu \bigcap_{k=1}^n \{X_{i(k)} \leq r_k\},$$

which verifies (6.3.4). \square

6.3.6 Embedding random variables in ZFC₋

Let L be an extralogical signature with $L_- \subseteq L$. If $P \subseteq \mathcal{L}^{\text{IS}}$ is an inductive theory with root $T_0 \supseteq \text{ZFC}_-$, then P is called a *real inductive theory in ZFC₋*. Note that, by definition, the root of an inductive theory is a consistent deductive theory. Hence, the existence of a real inductive theory in ZFC₋ presupposes the consistency of ZFC₋.

Let (S, Γ, ν) be a probability space and let $X = \langle X_i \mid i \in I \rangle$ be an indexed collection of real-valued random variables, with $\Gamma = \sigma(\langle X_i \mid i \in I \rangle)$. Let $C = \{\underline{X}_i \mid i \in I\}$ be a set of distinct constant symbols not in L_- , and define $L = L_- C$.

(T:prob-model-iso-ZFC-) **Theorem 6.3.4.** *Assume ZFC₋ is consistent. Then there exists a complete inductive theory $P \subseteq \mathcal{L}^{\text{IS}}$ with root ZFC₋ such that*

$$P(\bigwedge_{k=1}^n \underline{X}_{i(k)} \leq \underline{q}_k \mid \text{ZFC}_-) = \nu \bigcap_{k=1}^n \{X_{i(k)} \leq q_k\}, \quad (6.3.5) \quad \boxed{\text{prob-model-iso-ZFC-}}$$

whenever $i(1), \dots, i(n) \in I$ and $q_k \in \mathbb{Q}$.

Proof. We will call a set $J \subseteq \mathbb{R}$ a “rational interval” if it has one of the following five forms, for some $a, b \in \mathbb{Q}$: $J = \emptyset$, $J = (a, b]$, $J = (-\infty, b]$, $J = (a, \infty]$, $J = \mathbb{R}$. In \mathcal{L} , we adopt the following shorthand for every $a, b \in \mathbb{Q}$ with $a < b$:

- (i) $x \in (a, b] \leftrightarrow \underline{a} < x \wedge x < \underline{b} \vee x = \underline{b}$,
- (ii) $x \in (-\infty, b] \leftrightarrow x < \underline{b} \vee x = \underline{b}$, and

(iii) $x \in (a, \infty) \leftrightarrow a < x$.

In this way, if we adopt the shorthand $\underline{\emptyset} = \emptyset$, then we may write $x \in \underline{J}$ for every rational interval J .

A “rational cylinder” is a set $V \subseteq \mathbb{R}^I$ of the form

$$V = \{y \in \mathbb{R}^I \mid y_{i(1)} \in J_1, \dots, y_{i(n)} \in J_n\},$$

where $n \in \mathbb{N}$, $i(1), \dots, i(n) \in I$, and each J_k is a rational interval. If V is a rational cylinder, then we define the sentence $\varphi^V \in \mathcal{L}^0$ by

$$\varphi^V = \bigwedge_{k=1}^n \underline{X}_{i(k)} \in \underline{J}_k.$$

By Proposition 6.2.1, since ZFC_- is consistent, we may choose an L_- structure ω_0 such that $\omega_0 \models \Lambda_-^{\text{ZFC}}$. For each $y \in \mathbb{R}^I$, let $\omega = \omega^y$ be the L -expansion of ω_0 given by $\omega^{\underline{X}_i} = y_i$, and let $\Omega = \{\omega^y \mid y \in \mathbb{R}^I\}$. Let

$$\mathcal{E} = \{\varphi_\Omega^V \mid V \text{ is a rational cylinder}\},$$

and let Σ_0 be the set of finite, disjoint unions of sets in \mathcal{E} . Then Σ_0 is an algebra of sets on Ω .

Now define $\mathbb{P}_0 : \mathcal{E} \rightarrow [0, 1]$ by $\mathbb{P}_0 \varphi_\Omega^V = \nu \bigcap_{k=1}^n \{X_{i(k)} \in J_k\}$, and extend this to Σ_0 by finite additivity. Then \mathbb{P}_0 is a pre-measure on (Ω, Σ_0) with $\mathbb{P}_0 \Omega = 1$. By Carathéodory’s extension theorem (see, for instance, [10, Theorem 1.11]), there exists a unique probability measure \mathbb{P} on $(\Omega, \sigma(\Sigma_0))$ that agrees with \mathbb{P}_0 on Σ_0 .

Let $\Sigma = \sigma(\Sigma_0)$ and $\mathcal{P} = (\Omega, \Sigma, \mathbb{P})$. Since $\omega \models \text{ZFC}_-$ for all $\omega \in \Omega$, we have $\text{ZFC}_- \subseteq \text{Th } \mathcal{P}$. We may therefore define $P = \mathbf{Th } \mathcal{P} \downarrow_{[\text{ZFC}_-, \text{Th } \mathcal{P}]}$. By construction, (6.3.5) holds whenever each $q_k \in \mathbb{Q}$. \square

6.4 Real inductive theories in ZFC

(S:ind-th-ZFC)

In this section, we show how to represent real numbers in ZFC. We then use this representation to construct real inductive theories in ZFC. In stark contrast to what happens in ZFC_- , here we find that we can explicitly define every real number, every Borel set, and every measurable function. We will also construct a frame of reference in which each of these things is almost surely fixed, and not random. We call this the “real frame of reference,” and it is presented in three separate parts as Theorems 6.4.1, 6.4.3, and 6.4.5. Finally in Theorem 6.4.6, we show how real-valued random variables can be naturally and directly embedded into inductive models that are based on ZFC.

6.4.1 Real numbers and Borel sets

We now work in ZFC. We make all the same definitorial extensions that we did in ZFC_- , bringing us to the extralogical signature,

$$L = \{\in, \emptyset, \mathbf{s}, \underline{\mathbb{N}}_0, \underline{\mathbb{Z}}, \underline{\mathbb{Q}}, \underline{\mathbb{R}}, +, \cdot, <\} \cup \{q \mid q \in \mathbb{Q}\}.$$

This time, however, (6.3.2) holds, along with the analogous derivability relations for $\underline{\mathbb{Z}}$ and $\underline{\mathbb{Q}}$. In fact, for any $B \subseteq \mathbb{N}_0$, we can use AS to construct the set $\{x \in \underline{\mathbb{N}}_0 \mid \bigvee_{n \in B} x = \underline{n}\}$. In other words, we can explicitly define every set in $\mathfrak{P} \underline{\mathbb{N}}_0$, and for each of them, the analogue of (6.3.2) holds. The same is true for every set of integers and every set of rationals.

Since $\underline{\mathbb{R}}$ is defined as a set of Dedekind cuts, and a Dedekind cut is a set of rationals, it follows that we can explicitly define each individual real number. We may therefore add, for each $r \in \mathbb{R}$, an explicitly defined constant symbol \underline{r} , using the infinitary definition,

$$\delta_r(y) = \forall x(x \in y \leftrightarrow \bigvee_{q \in B(r)} x = \underline{q}),$$

where $B(r) = \{q \in \mathbb{Q} \mid q < r\}$.

We now add to ZFC an explicit, finitary definition of $\underline{\mathcal{B}}$, which denotes the Borel σ -algebra, $\mathcal{B}(\mathbb{R})$, on \mathbb{R} . Unlike with \mathbb{N}_0 , \mathbb{Z} , and \mathbb{Q} , we do not expect to be able to explicitly define each individual subset of \mathbb{R} , since we cannot form an uncountable disjunction in the language \mathcal{L} . We can, however, explicitly define each Borel set.

Let $\mathcal{E} \subseteq \mathfrak{P} \mathbb{R}$ be given by $\mathcal{E} = \{(-\infty, r] \mid r \in \mathbb{R}\}$. Note that $\mathcal{B}(\mathbb{R}) = \sigma(\mathcal{E})$. Recall the recursive construction of $\mathcal{B}(\mathbb{R})$ from \mathcal{E} given in Section 2.3.1, and recall that if $V \in \mathcal{B}(\mathbb{R})$, then $\text{rk } V$ denotes the rank of V with respect to \mathcal{E} . For each $V \in \mathcal{B}(\mathbb{R})$, we explicitly define \underline{V} by recursion on $\text{rk } V$.

If $\text{rk } V = 0$, then $V \in \mathcal{E}$, so that $V = (-\infty, r]$. Define $\underline{V} = \{x \in \underline{\mathbb{R}} \mid x \leq \underline{r}\}$, where we adopt the shorthand $x \leq y$ for $x < y \vee x = y$. If $\text{rk } V = \alpha$, then α is a successor ordinal, and we may write $\alpha = \beta + 1$ for some β . If $V \in \mathcal{E}'_\beta$, then we may choose $W \in \mathcal{E}_\beta$ such that $V = W^c$. We then define $\underline{V} = \{x \in \underline{\mathbb{R}} \mid x \notin \underline{W}\}$. If $V \notin \mathcal{E}'_\beta$, then we may choose a nonempty and countable $\mathcal{D} \subseteq \mathcal{E}'_\beta$ such that $V = \bigcap \mathcal{D}$. We then define $\underline{V} = \{x \in \underline{\mathbb{R}} \mid \bigwedge_{W \in \mathcal{D}} x \in \underline{W}\}$. This defines \underline{V} for every $V \in \mathcal{B}(\mathbb{R})$.

We now have the extralogical signature,

$$L_{\text{ZFC}} = \{\in, \emptyset, \mathbf{s}, \underline{\mathbb{N}}_0, \underline{\mathbb{Z}}, \underline{\mathbb{Q}}, \underline{\mathbb{R}}, \underline{\mathcal{B}}, +, \cdot, <\} \cup \{\underline{r} \mid r \in \mathbb{R}\} \cup \{\underline{V} \mid V \in \mathcal{B}(\mathbb{R})\}.$$

(6.4.1) ZFC-signature

As usual, we have omitted the duplicates, and each symbol in L_{ZFC} , except for \in , is an explicitly defined constant symbol.

6.4.2 The real frame of reference

Let L be an extralogical signature that contains a binary relation symbol \in . Let $\omega = (A, L^\omega)$ be an L -structure. For a given $b \in A$, we define $\omega b = \{a \in A \mid a \in^\omega b\}$. If t is a ground term, then we write $\omega t = \omega(t^\omega)$.

(T:real-FOR)

Theorem 6.4.1 (Real frame of reference I). *Let L be an extralogical signature such that $L_{\text{ZFC}} \subseteq L$. Let \mathcal{Q} be an L -model such that $\mathcal{Q} \models \text{ZFC}$. Then there exists an L -model $\mathcal{P} = (\Omega, \Sigma, \mathbb{P})$ such that $\mathcal{Q} \simeq \mathcal{P}$ and*

- (i) $\omega \models \text{ZFC}_-$ for every $\omega \in \Omega$,

- (ii) $\underline{q}^\omega = q$ for every $\omega \in \Omega$,
- (iii) ${}^\omega\mathbb{R} \subseteq \mathbb{R}$ for every $\omega \in \Omega$,
- (iv) $\underline{r}^\omega = r$ a.s., for each $r \in \mathbb{R}$, and
- (v) $\underline{V}^\omega = {}^\omega\underline{V} = V \cap {}^\omega\mathbb{R}$ a.s., for each $V \in \mathcal{B}(\mathbb{R})$.

Proof. Let $\mathcal{Q} = (\Omega', \Sigma', \mathbb{Q})$ be an L -model such that $\mathcal{Q} \models \text{ZFC}$. Let

$$L'_{\text{ZFC}} = \{\emptyset, \mathfrak{S}, \underline{\mathbb{N}_0}, \underline{\mathbb{Z}}, \underline{\mathbb{Q}}, \underline{\mathbb{R}}, \underline{\mathcal{B}}, +, \cdot, <\} \cup \{\underline{q} \mid q \in \mathbb{Q}\},$$

so that L'_{ZFC} is countable and every $\mathfrak{s} \in L'_{\text{ZFC}}$ is an explicitly defined constant symbol with defining formula $\delta_{\mathfrak{s}}(y)$. Let

$$\varphi = \forall y (y \in \underline{\mathbb{Q}} \leftrightarrow \bigvee_{q \in \mathbb{Q}} y = \underline{q}),$$

and note that $\text{ZFC} \vdash \varphi$. Let

$$X_0 = \Lambda_{\text{ZFC}} \cup \{\varphi\} \cup \{\exists! y \delta_{\mathfrak{s}}(y) \mid \mathfrak{s} \in L'_{\text{ZFC}}\},$$

so that $\mathcal{Q} \models X_0$. Since X_0 is countable, it follows that $\nu \models X_0$ for \mathbb{Q} -a.e. $\nu \in \Omega'$. By Remark 5.3.9, we may assume that $\nu \models X_0$ for all $\nu \in \Omega'$.

Fix $\nu = (A_\nu, L^\nu) \in \Omega'$. Let $a \in {}^\nu\mathbb{R}$ and define $B = \{q \in \mathbb{Q} \mid \nu \models (q < x)[a]\}$. Since $\nu \models (\forall x \in \mathbb{R})(\exists y \in \mathbb{Q})(y < x)$ and $\nu \models \varphi$, it follows that $B \neq \emptyset$. Similarly, $B \neq \mathbb{Q}$, B is downward closed, and B has no maximum element. That is, B is a Dedekind cut. Therefore, there exists a unique $r \in \mathbb{R}$ such that $B = B(r) = \{q \in \mathbb{Q} \mid q < r\}$. Let $g_\nu : {}^\nu\mathbb{R} \rightarrow \mathbb{R}$ denote the function $a \mapsto r$.

Let $a, a' \in {}^\nu\mathbb{R}$ and define B and B' accordingly. Assume $a \neq a'$. Since

$$\nu \models (\forall xy \in \mathbb{R})(x \neq y \rightarrow x < y \vee y < x),$$

we have $a <^\nu a'$ or $a' <^\nu a$. Without loss of generality, assume it is the former. Since

$$\nu \models (\forall xy \in \mathbb{R})(x < y \rightarrow (\exists z \in \mathbb{Q})(x < z \wedge z < y)),$$

and $\nu \models \varphi$, there exists $q \in \mathbb{Q}$ such that $a <^\nu \underline{q}$ and $\underline{q} <^\nu a'$. We then have $q \in B' \setminus B$, so that $B \neq B'$, which implies $g_\nu a \neq g_\nu a'$. Therefore, g_ν is injective. In fact, g_ν is order-preserving, in the sense that $a <^\nu a'$ if and only if $g_\nu a < g_\nu a'$.

Now let $b \in {}^\nu(\mathfrak{P}\mathbb{R})$ and $a \in {}^\nu b$. Then $\nu \models (x \in y \wedge y \in \mathfrak{P}\mathbb{R})[a, b]$. Hence, we have $\nu \models (x \in \mathbb{R})[a]$, so that $a \in {}^\nu\mathbb{R}$. This shows that ${}^\nu b \subseteq {}^\nu\mathbb{R}$, so that $g_\nu {}^\nu b \subseteq \mathbb{R}$. Since ${}^\nu\mathbb{R} \cap {}^\nu(\mathfrak{P}\mathbb{R}) = \emptyset$, we may extend g_ν to be a function $g_\nu : {}^\nu\mathbb{R} \cup {}^\nu(\mathfrak{P}\mathbb{R}) \rightarrow \mathbb{R} \cup \mathfrak{P}\mathbb{R}$ by setting $g_\nu b = g_\nu {}^\nu b$ for all $b \in {}^\nu(\mathfrak{P}\mathbb{R})$. As above, it follows that g_ν remains injective.

We now extend g_ν to be a bijection from A_ν to some set A_ω , which is a superset of $g_\nu({}^\nu\mathbb{R} \cup {}^\nu(\mathfrak{P}\mathbb{R}))$. Let ω be the isomorphic image of ν under g_ν . Note that since g_ν is order-preserving on ${}^\nu\mathbb{R}$, we have $a <^\omega a'$ if and only if $a < a'$, whenever $a, a' \in A_\omega \cap \mathbb{R}$.

Let h denote the function $\nu \mapsto \omega$ and let $\Omega = h\Omega'$. Let $\mathcal{P} = (\Omega, \Sigma, \mathbb{P})$ be the measure space image of \mathcal{Q} under h . Then h induces a measure-space

isomorphism from \mathcal{Q} to \mathcal{P} , and $\nu \simeq h\nu$ for all $\nu \in \Omega'$. Hence, h is a model isomorphism, and $\mathcal{Q} \simeq \mathcal{P}$. Since $\nu \models X_0$ for all $\nu \in \Omega'$ and $\omega \simeq \nu$, we have (i). Moreover, $a \in {}^\omega\mathbb{R}$ if and only if $g_\nu^{-1}a \in {}^\nu\mathbb{R}$. Hence, ${}^\omega\mathbb{R} = g_\nu {}^\nu\mathbb{R} \subseteq \mathbb{R}$, so (iii) holds.

Now let $r \in \mathbb{R}$. Since $\mathcal{P} \models \exists!y \delta_r(y)$, we may choose $\Omega^* \in \Sigma$ such that $\mathbb{P}\Omega^* = 1$ and $\omega \models \exists!y \delta_r(y)$ for all $\omega \in \Omega^*$. Fix $\omega \in \Omega^*$. Then $\omega \models \underline{r} \in \mathbb{R}$. Hence, $\underline{r}^\omega \in {}^\omega\mathbb{R} = g_\nu {}^\nu\mathbb{R}$, so that

$$\begin{aligned} \underline{r}^\omega &= g_\nu g_\nu^{-1} \underline{r}^\omega = \sup\{q \in \mathbb{Q} \mid \nu \models (q < x)[g_\nu^{-1} \underline{r}^\omega]\} \\ &= \sup\{q \in \mathbb{Q} \mid \omega \models (q < x)[\underline{r}^\omega]\} \\ &= \sup\{q \in \mathbb{Q} \mid \omega \models \underline{q} < r\}. \end{aligned}$$

But by the definition of \underline{r} , we have

$$\omega \models (\forall x \in \mathbb{Q})(x < \underline{r} \leftrightarrow \bigvee_{q \in B(r)} x = q),$$

where $B(r) = \{q \in \mathbb{Q} \mid q < r\}$. Thus, $\omega \models \underline{q} < \underline{r}$ if and only if $q \in B(r)$. Therefore, $\underline{r}^\omega = \sup B(r) = r$, proving (iv).

Now fix $V \in \mathcal{B}(\mathbb{R})$. We will prove (v) by induction on $\text{rk } V$. Suppose $\text{rk } V = 0$. Then $V = (-\infty, r]$. Choose $\Omega^* \in \Sigma$ such that $\mathbb{P}\Omega^* = 1$ and, for all $\omega \in \Omega^*$, we have $\underline{r}^\omega = r$ and $\omega \models \exists!y \delta_V(y)$, where $\delta_V(y)$ is the defining formula for \underline{V} . Since $\omega \models \underline{V} \in \mathfrak{P}\mathbb{R}$, we have $\nu \models \underline{V} \in \mathfrak{P}\mathbb{R}$, so that $\underline{V}^\nu \in {}^\nu(\mathfrak{P}\mathbb{R})$. Therefore, $\underline{V}^\omega = g_\nu {}^\nu \underline{V}^\nu = g_\nu {}^\nu \underline{V}$. On the other hand,

$$\begin{aligned} a \in {}^\omega \underline{V} &\text{ iff } a \in A_\omega \text{ and } \omega \models (x \in \underline{V})[a] \\ &\text{ iff } g_\nu^{-1}a \in A_\nu \text{ and } \nu \models (x \in \underline{V})[g_\nu^{-1}a] \\ &\text{ iff } g_\nu^{-1}a \in {}^\nu \underline{V}. \end{aligned}$$

Thus, ${}^\omega \underline{V} = g_\nu {}^\nu \underline{V} = \underline{V}^\omega$. Now, $\underline{V}^\omega = \{x \in \mathbb{R} \mid x \leq \underline{r}\}^\omega$. Hence,

$${}^\omega \underline{V} = \{a \in {}^\omega \mathbb{R} \mid \omega \models (x \leq \underline{r})[a]\} = \{a \in {}^\omega \mathbb{R} \mid a \leq^\omega \underline{r}^\omega\}.$$

But ${}^\omega \mathbb{R} \subseteq \mathbb{R}$ and $\underline{r}^\omega = r \in \mathbb{R}$. Therefore, $a \leq^\omega \underline{r}^\omega$ if and only if $a \leq r$, which implies $\underline{V}^\omega = {}^\omega \underline{V} = V \cap {}^\omega \mathbb{R}$.

Now suppose $\text{rk } V = \alpha$. Since the rank of a Borel set is always a successor ordinal, we may write $\alpha = \beta + 1$. Assume $V \in \mathcal{E}'_\beta$. Then $\underline{V} = \{x \in \mathbb{R} \mid x \notin \underline{W}\}$, where $W \in \mathcal{E}_\beta$ and $V = \mathbb{R} \setminus W$. As above, choose $\Omega^* \in \Sigma$ such that $\mathbb{P}\Omega^* = 1$ and, for all $\omega \in \Omega^*$, we have $\omega \models \exists!y \delta_W(y)$ and $\omega \models \exists!y \delta_V(y)$. Then

$$\underline{V}^\omega = {}^\omega \underline{V} = \{a \in {}^\omega \mathbb{R} \mid a \notin {}^\omega \underline{W}\} = {}^\omega \mathbb{R} \setminus (W \cap {}^\omega \mathbb{R}) = V \cap {}^\omega \mathbb{R},$$

proving (v) in the case that $V \in \mathcal{E}'_\beta$. The proof in the case $V \notin \mathcal{E}'_\beta$ is similar.

Finally, since $\underline{q}^\omega = q$ a.s., for each $q \in \mathbb{Q}$, and \mathbb{Q} is countable, it follows that $\underline{q}^\omega = q$ for each $q \in \mathbb{Q}$, a.s. By Remark 5.3.9, we may take \mathcal{P} to be such that (ii) holds. \square

(C:real-FOR) **Corollary 6.4.2.** *Let $r \in \mathbb{R}$ and $V, V', V_n \in \mathcal{B}(\mathbb{R})$.*

- (i) If $r \in V$, then $\text{ZFC} \vdash \underline{r} \in \underline{V}$.
- (ii) If $r \notin V$, then $\text{ZFC} \vdash \underline{r} \notin \underline{V}$.
- (iii) If $V \subseteq V'$, then $\text{ZFC} \vdash \underline{V} \subseteq \underline{V}'$.
- (iv) If $V = \bigcap_{n \in \mathbb{N}_0} V_n$, then $\text{ZFC} \vdash \forall x (x \in \underline{V} \leftrightarrow \bigvee_{n \in \mathbb{N}_0} x \in \underline{V}_n)$.

Proof. Let $r \in V$. Suppose $\mathcal{P} \models \text{ZFC}$. Using the real frame of reference, we may assume \mathcal{P} satisfies (i)–(v) in Theorem 6.4.1. Hence, for a.e. ω , we have $\underline{r}^\omega = r$. Since $\underline{r}^\omega \in {}^\omega \underline{\mathbb{R}}$, this gives $\underline{r}^\omega \in V \cap {}^\omega \underline{\mathbb{R}} = {}^\omega \underline{V}$, so that $\omega \models (\underline{r} \in \underline{V})$. Since this holds almost surely, we have $\mathcal{P} \models \underline{r} \in \underline{V}$. The proofs of (ii)–(iv) are similar. \square

6.4.3 Sequences and limits

Let us define, in ZFC, subtraction of real numbers and the absolute value function, and adopt the usual shorthand for these functions. Let

$$\begin{aligned} \varphi_{\text{lim}}(u, v) = & u \in \underline{\mathbb{R}}^{\mathbb{N}_0} \wedge v \in \underline{\mathbb{R}} \\ & \wedge (\forall z \in (0, \infty)) (\exists y \in \underline{\mathbb{N}}_0) (\forall x \in \underline{\mathbb{N}}_0) (x \geq y \rightarrow |u(x) - v| < z). \end{aligned}$$

Then $\varphi_{\text{lim}}(u, v)$ says that u is a sequence of real numbers that converges to the real number v .

(T:real-FOR-2) **Theorem 6.4.3 (Real frame of reference II).** *The model $\mathcal{P} = (\Omega, \Sigma, \mathbb{P})$ in Theorem 6.4.1 may be chosen so that for every $\omega \in \Omega$, we have the following:*

- (i) for each $a \in {}^\omega(\underline{\mathbb{R}}^{\mathbb{N}_0})$ and each $n \in \mathbb{N}_0$, there exists a unique $a_n \in {}^\omega \underline{\mathbb{R}} \subseteq \mathbb{R}$ such that $\omega \models (y = x(\underline{n}))[a, a_n]$, and
- (ii) $\omega \models \varphi_{\text{lim}}[a, b]$ if and only if $b \in {}^\omega \underline{\mathbb{R}}$ and $a_n \rightarrow b$.

Proof. Note that $\text{ZFC} \vdash \psi$, where

$$\psi = x \in \underline{\mathbb{R}}^{\mathbb{N}_0} \rightarrow \bigwedge_{n \in \mathbb{N}_0} (\exists! y \in \underline{\mathbb{R}}) y = x(\underline{n}).$$

Hence, $\omega \models \psi$ a.s. By Remark 5.3.9, we may assume $\omega \models \psi$ for all $\omega \in \Omega$. Fix $\omega \in \Omega$ and suppose $a \in {}^\omega(\underline{\mathbb{R}}^{\mathbb{N}_0})$. Then $\omega \models (x \in \underline{\mathbb{R}}^{\mathbb{N}_0})[a]$. Hence, since $\omega \models \psi$, it follows that for each $n \in \mathbb{N}_0$, there exists a unique $a_n \in {}^\omega \underline{\mathbb{R}} \subseteq \mathbb{R}$ such that $\omega \models (y = x(\underline{n}))[a, a_n]$. Note that in the definition of φ_{lim} , it suffices to consider rational z . Since $\underline{q}^\omega = q$ for all $q \in \mathbb{Q}$, it follows that $\omega \models \varphi_{\text{lim}}[a, b]$ if and only if $b \in {}^\omega \underline{\mathbb{R}}$ and $a_n \rightarrow b$. \square

Remark 6.4.4. Let $\mathcal{P} = (\Omega, \Sigma, \mathbb{P})$ be a model that satisfies Theorem 6.4.3. In \mathcal{L} , we can formulate a sentence φ which asserts that $+$ is continuous. For such a sentence, $\text{ZFC}_- \vdash \varphi$. Recall that $+$ was defined so that it agrees with the usual addition of rational numbers. That is, $\text{ZFC}_- \vdash \underline{q}_1 + \underline{q}_2 = \underline{q}_3$ if and only if $q_1 + q_2 = q_3$, whenever $q_i \in \mathbb{Q}$. Using this, it can be shown that, for all $\omega \in \Omega$ and all $a, b, c \in {}^\omega \underline{\mathbb{R}} \subseteq \mathbb{R}$, we have $\omega \models (x + y = z)[a, b, c]$ if and only if $a + b = c$. A similar thing holds for \cdot and $<$.

6.4.4 Measurable functions

Let $h : \mathbb{R} \rightarrow \mathbb{R}$ be measurable. Suppose h is a simple function. That is, the range of h is a finite set $\{r_1, \dots, r_n\}$. Let $V_j = h^{-1}\{r_j\}$. Then we may explicitly define h in ZFC by

$$\underline{h} = \{(x, y) \in \mathbb{R} \times \mathbb{R} \mid \bigvee_{i=1}^n x \in \underline{V}_i \wedge y = r_i\}.$$

Suppose h is not a simple function. Choose a sequence $\langle h_n \mid n \in \mathbb{N}_0 \rangle$ of simple functions such that $h_n \rightarrow h$ pointwise. In ZFC, we define

$$\underline{h}_. = \{(x, y) \in \mathbb{R} \times \mathbb{R}^{\mathbb{N}_0} \mid \bigwedge_{n \in \mathbb{N}_0} y(n) = \underline{h}_.(x)\}.$$

Then $\underline{h}_.$ is an explicit definition of the function $x \mapsto \langle h_n(x) \mid n \in \mathbb{N}_0 \rangle$. We now define

$$\underline{h} = \{(x, y) \in \mathbb{R} \times \mathbb{R} \mid \varphi_{\lim}(\underline{h}_.(x), y)\}.$$

(T:real-FOR-3) **Theorem 6.4.5 (Real frame of reference III).** *Let $h : \mathbb{R} \rightarrow \mathbb{R}$ be a measurable function and let $\mathcal{P} = (\Omega, \Sigma, \mathbb{P})$ be a model chosen according to Theorem 6.4.3. Then, for \mathbb{P} -a.e. $\omega \in \Omega$, we have $\omega \models (y = \underline{h}(x))[a, b]$ if and only if $a \in {}^\omega\mathbb{R}$, $b \in {}^\omega\mathbb{R}$, and $h(a) = b$.*

Proof. First suppose h is a simple function. Then $\omega \models (y = \underline{h}(x))[a, b]$ if and only if there exists $i \in \{1, \dots, n\}$ such that $a \in {}^\omega V_i$ and $b = r_i^\omega$. By Theorem 6.4.1, we may choose $\Omega^* \in \Sigma$ such that $\mathbb{P}\Omega^* = 1$ and, for all $\omega \in \Omega^*$, we have ${}^\omega\mathbb{R} \subseteq \mathbb{R}$, $r_i^\omega = r_i$, and ${}^\omega V_i = V_i \cap {}^\omega\mathbb{R}$. Hence, the conclusion of the theorem holds for each $\omega \in \Omega^*$.

Now suppose h is not a simple function. Let $\langle h_n \mid n \in \mathbb{N}_0 \rangle$ be the sequence of simple functions used to define $\underline{h}_.$. Choose $\Omega^* \in \Sigma$ such that $\mathbb{P}\Omega^* = 1$ and $\omega \models \exists! y \delta_s(y)$ for every $s \in \{\underline{h}_. \mid n \in \mathbb{N}_0\} \cup \{\underline{h}_., \underline{h}_.\}$. Suppose $\omega \models (y = \underline{h}(x))[a, b]$. Then $\omega \models \varphi_{\lim}(\underline{h}_.(x), y)[a, b]$, which means $\omega \models \varphi_{\lim}[c, b]$, where $\omega \models (u = \underline{h}_.(x))[a, c]$. By Theorem 6.4.3(i), for each $n \in \mathbb{N}_0$, there exists a unique $c_n \in {}^\omega\mathbb{R} \subseteq \mathbb{R}$ such that $\omega \models (v = u(n))[c, c_n]$. Since $\omega \models \varphi_{\lim}[c, b]$, Theorem 6.4.3(ii) implies $b \in {}^\omega\mathbb{R}$ and $c_n \rightarrow b$. On the other hand, since $\omega \models (u = \underline{h}_.(x))[a, c]$, it follows from the definition of $\underline{h}_.$ that $\omega \models (u(n) = \underline{h}_.(x))[a, c]$. Therefore, $\omega \models (v = \underline{h}_.(x))[a, c_n]$. Since the theorem holds for simple functions, we have $c_n = h_n(a)$. Thus, $b = \lim h_n(a) = h(a)$. The proof of the converse is similar. \square

6.4.5 Embedding random variables in ZFC

Let L be an extralogical signature with $L_{\text{ZFC}} \subseteq L$. If $P \subseteq \mathcal{L}^{\text{IS}}$ is an inductive theory with root $T_0 \supseteq \text{ZFC}$, then P is called a *real inductive theory in ZFC*. As with ZFC_- , the existence of a real inductive theory in ZFC presupposes the consistency of ZFC.

Let (S, Γ, ν) be a probability space and let $X = \langle X_i \mid i \in I \rangle$ be an indexed collection of real-valued random variables, with $\Gamma = \sigma(\langle X_i \mid i \in I \rangle)$. Let

$C = \{\underline{X}_i \mid i \in I\}$ be a set of distinct constant symbols not in L_{ZFC} , and define $L = L_{\text{ZFC}}C$.

Theorem 6.4.6 below shows that, if ZFC is strictly satisfiable, then we recover the full analogue of Theorem 5.4.2, including the natural correspondence between outcomes and structures. Recall from Theorem 6.2.4 that the existence of a strongly inaccessible cardinal implies not only the consistency of ZFC_- , but also the strict satisfiability of ZFC. Hence, the assumption in Theorem 6.4.6 is not too far beyond the typically uncontroversial assumption that ZFC_- is consistent.

The proof of Theorem 6.4.6 shows how the inductive model $\mathcal{P} = (\Omega, \Sigma, \mathbb{P})$ is built from the measure-theoretic probability model (S, Γ, ν, X) . Each structure $\omega \in \Omega$ is an expansion of a single structure ω_0 that strictly satisfies ZFC. This implies, for example, that \underline{r}^ω does not depend on ω . The same is true for every other symbol in L_{ZFC} . In other words, all the familiar objects of set theory and the real numbers are all fixed in \mathcal{P} . The only things that vary with ω are \underline{X}_i^ω . This exactly matches our intuition about measure-theoretic models. In practice, when we work with a measure-theoretic model, we think of the random variables as being the only things that are “random.” The real numbers, their relations, the concept of set membership, and so on, are all fixed, and do not change from one outcome to another. In a sense, then, the practicing probabilist, in using the logic of countable unions and intersections, and in assuming that all our familiar mathematics is fixed under different outcomes, is operating under the implicit assumption that ZFC is strictly satisfiable.

$\langle \text{T:prob-model-iso-ZFC} \rangle$ **Theorem 6.4.6.** *Assume ZFC is strictly satisfiable. Then there exists an \mathcal{L} -model $\mathcal{P} = (\Omega, \Sigma, \mathbb{P})$ with $\mathcal{P} \models \text{ZFC}$, and a function $h : S \rightarrow \Omega$ mapping $x \in S$ to $\omega \in \Omega$ such that*

- (i) $x \in \{X_i \in V\}$ if and only if $\omega \models \underline{X}_i \in \underline{V}$ for all $V \in \mathcal{B}(\mathbb{R})$,
- (ii) each $U \in \Gamma$ can be written as $U = h^{-1}\varphi_\Omega$ for some $\varphi \in \mathcal{L}^0$, and
- (iii) h induces a measure-space isomorphism from (S, Γ, ν) to \mathcal{P} .

Consequently, if $P = \mathbf{Th} \mathcal{P} \upharpoonright_{[\text{ZFC}, \text{Th} \mathcal{P}]}$, then P satisfies

$$P(\bigwedge_{k=1}^n \underline{X}_{i(k)} \in \underline{V}_k \mid \text{ZFC}) = \nu \bigcap_{k=1}^n \{X_{i(k)} \in V_k\}, \quad (6.4.2) \quad \boxed{\text{prob-model-iso-ZFC}}$$

whenever $i(1), \dots, i(n) \in I$ and $V_k \in \mathcal{B}(\mathbb{R})$.

Proof. Assume there exists an L_{ZFC} -structure ω_0 such that $\omega_0 \models \text{ZFC}$. For each $x \in S$, define $\omega = \omega^x$ to be the L -expansion of ω_0 given by $\underline{X}_i^\omega = \underline{r}^{\omega_0}$, where $r = X_i(x)$. Let $\Omega = \{\omega^x \mid x \in S\}$ and let $h : S \rightarrow \Omega$ denote the map $x \mapsto \omega^x$. Let $\mathcal{P} = (\Omega, \Sigma, \mathbb{P})$ be the measure space image of (S, Γ, ν) under h . Since $\omega \models \text{ZFC}$ for all $\omega \in \Omega$, we have $\mathcal{P} \models \text{ZFC}$.

Let $x \in S$ and $V \in \mathcal{B}(\mathbb{R})$, and set $r = X_i(x)$. By the definition of ω^x , we have $\omega^x \models \underline{X}_i \in \underline{V}$ if and only if $\omega_0 \models \underline{r} \in \underline{V}$. By Corollary 6.4.2, we have $\text{ZFC} \vdash \underline{r} \in \underline{V}$ if $r \in V$, and $\text{ZFC} \vdash \underline{r} \notin \underline{V}$ if $r \notin V$. Hence, since $\omega_0 \models \text{ZFC}$, we

have $\omega_0 \models r \in \underline{V}$ if and only if $r \in V$. But $r = X_i(x)$. Hence, $x \in \{X_i \in V\}$ if and only if $\omega^x \models \underline{X}_i \in \underline{V}$, proving (i). It follows as in the proof of Theorem 6.3.3 that (ii) and (iii) hold, and (6.4.2) holds for $P = \mathbf{Th} \mathcal{P} \downarrow_{[\mathbf{ZFC}, Th \mathcal{P}]}$. \square

6.5 Limit theorems

(S:limit-thms)

In this section, we introduce the notion of Borel terms, and use them to formulate, in inductive logic, both the law of large numbers and the central limit theorem.

6.5.1 Borel terms

Let P be a real inductive theory in ZFC. That is, $P \subseteq \mathcal{L}^{\text{IS}}$, where $L_{\text{ZFC}} \subseteq L$, and P has root $T_0 \supseteq \text{ZFC}$. In particular, we are assuming ZFC is consistent.

Let $X \in \text{ante } P$. A ground term $t \in \mathcal{T}$ is *real given* X if $P(t \in \underline{\mathbb{R}} \mid X) = 1$. We say that t is *Borel given* X if it is real given X and $P(t \in \underline{V} \mid X)$ exists for all $V \in \mathcal{B}(\mathbb{R})$. If t is Borel given X , then the *distribution of t given X* is the function $\mu = \mu_{t|X}$ from $\mathcal{B}(\mathbb{R})$ to $[0, 1]$ given by $\mu V = P(t \in \underline{V} \mid X)$.

(P:term-dist)

Proposition 6.5.1. *If t is Borel given X , then the distribution of t is a probability measure on $(\mathbb{R}, \mathcal{B}(\mathbb{R}))$.*

Proof. Let $\mathcal{Q} = (\Omega, \Gamma, \mathbb{Q}) \models P$. We may write $X \equiv Y \cup \{\psi\}$, where $\mathcal{Q} \models Y$ and $P(\varphi \mid X) = \overline{\mathbb{Q}} \varphi_{\Omega} \cap \psi_{\Omega} / \overline{\mathbb{Q}} \psi_{\Omega}$ whenever $P(\varphi \mid X)$ is defined. Let $\Sigma = \overline{\Gamma}$ and let \mathbb{P} be the probability measure on (Ω, Σ) given by $\mathbb{P} A = \overline{\mathbb{Q}} A \cap \psi_{\Omega} / \overline{\mathbb{Q}} \psi_{\Omega}$. Then $P(\varphi \mid X) = \mathbb{P} \varphi_{\Omega}$ whenever $P(\varphi \mid X)$ is defined. Since $\mathcal{P} = (\Omega, \Sigma, \mathbb{P}) \models \text{ZFC}$, we may use the real frame of reference, and assume \mathcal{P} satisfies (i)–(v) in Theorem 6.4.1.

Thus, ${}^{\omega} \emptyset = \emptyset$ a.s. Hence, for a.e. ω , we have $t^{\omega} \notin {}^{\omega} \emptyset$, meaning $\omega \not\models t \in \emptyset$. This gives $(t \in \emptyset)_{\Omega} = \emptyset$ a.s., so that $\mu \emptyset = P(t \in \emptyset \mid X) = 0$. Similarly, since $P(t \in \underline{\mathbb{R}} \mid X) = 1$, we have $\omega \models t \in \underline{\mathbb{R}}$ a.s., and it follows that $\mu \mathbb{R} = 1$.

Let $\{V_n\} \subseteq \mathcal{B}(\mathbb{R})$ be pairwise disjoint, and set $V = \bigcup_n V_n$. Define the sentences $\varphi_n = t \in \underline{V}_n$. Let $i \neq j$. As in the proof of Corollary 6.4.2, we have $\text{ZFC} \vdash \forall x \neg(x \in \underline{V}_i \wedge x \in \underline{V}_j)$. Hence, $\text{ZFC} \vdash \neg(\varphi_i \wedge \varphi_j)$, which gives $P(\varphi_i \wedge \varphi_j \mid X) = 0$. Theorem 3.2.24 therefore implies $P(\bigvee_n \varphi_n \mid X) = \sum_n P(\varphi_n \mid X) = \sum_n \mu V_n$. On the other hand, as above, $\text{ZFC} \vdash t \in \underline{V} \leftrightarrow \bigvee_n \varphi_n$, so that Proposition 3.2.14 gives $P(\bigvee_n \varphi_n \mid X) = P(t \in V \mid X) = \mu V$. \square

Let t be Borel given X . We define the *expected value of t given X* by $E[t \mid X] = \int_{\mathbb{R}} x \mu_{t|X}(dx)$, provided this integral exists. Note that $E[t \mid X] \in [-\infty, \infty]$. If $\int_{\mathbb{R}} |x| \mu_{t|X}(dx) < \infty$, then we say t is *integrable given X* . In this case, $E[t \mid X] \in \mathbb{R}$.

6.5.2 Jointly Borel terms

A finite sequence of terms, t_1, \dots, t_n , are *jointly Borel given X* if each t_i is Borel given X and $P(\bigwedge_{i=1}^n t_i \in \underline{V}_i \mid X)$ exists, whenever $V_i \in \mathcal{B}(\mathbb{R})$. If

t_1, \dots, t_n are jointly Borel, then each t_i is Borel. The converse holds when P is complete, by Definition 3.3.1(i). More generally, an indexed collection of terms, $\mathbf{t} = \langle t_i \mid i \in I \rangle$, is *jointly Borel given X* if each finite subsequence is.

$\langle \text{P:term-dist-Rn} \rangle$ **Proposition 6.5.2.** *If $\mathbf{t} = (t_1, \dots, t_n)$ are jointly Borel given X , then there exists a unique probability measure $\mu = \mu_{\mathbf{t}|X}$ on $(\mathbb{R}^n, \mathcal{B}(\mathbb{R}^n))$ such that*

$$\mu V_1 \times \cdots \times V_n = P(\bigwedge_{k=1}^n t \in \underline{V}_k \mid X),$$

whenever $i(1), \dots, i(n) \in I$ and $V_k \in \mathcal{B}(\mathbb{R})$.

Proof. The existence and uniqueness of μ follows from Carathéodory's extension theorem, as in the proof of Theorem 6.3.4, using as our algebra the set of finite disjoint unions of measurable rectangles, $V_1 \times \cdots \times V_n$. \square

The unique Borel probability measure $\mu_{\mathbf{t}|X}$ in Proposition 6.5.2 is called the *distribution of \mathbf{t} given X* .

$\langle \text{R:term-dist-Rn} \rangle$ **Remark 6.5.3.** Just as we did for $\mathcal{B}(\mathbb{R})$, we can explicitly define \underline{V} for each $V \in \mathcal{B}(\mathbb{R}^n)$ and establish a real frame of reference for \mathbb{R}^n . If μ is the distribution of (t_1, \dots, t_n) given X , then $\mu V = P((t_1, \dots, t_n) \in \underline{V} \mid X)$ for all $V \in \mathcal{B}(\mathbb{R}^n)$. To see this, let Γ be the set of $V \subseteq \mathbb{R}^n$ for which it holds. Proposition 6.5.2 shows that Γ contains the algebra of finite disjoint unions of measurable rectangles, which is a π -system that generates $\mathcal{B}(\mathbb{R}^n)$. Since P is an inductive theory, we have that Γ is a λ -system. Therefore, by Dynkin's π - λ theorem, $\mathcal{B}(\mathbb{R}^n) \subseteq \Gamma$.

6.5.3 Independence of terms

Let I be a set with $|I| \geq 2$ and let $\langle t_i \mid i \in I \rangle$ be jointly Borel given X . Then $\langle t_i \mid i \in I \rangle$ are *independent given X* if $\langle t_i \in \underline{V}_i \mid i \in I \rangle$ are independent given X , whenever $V_i \in \mathcal{B}(\mathbb{R})$. We say that $\langle t_i \mid i \in I \rangle$ are *identically distributed given X* if the distribution of t_i given X does not depend on i . As usual, we write i.i.d. for the phrase, "independent and identically distributed."

$\langle \text{R:dist-indep} \rangle$ **Remark 6.5.4.** If $\mathbf{t} = (t_1, \dots, t_n)$ are independent given X , then Corollary 4.5.12 implies that $\mu_{\mathbf{t}|X} = \prod_{k=1}^n \mu_{t_k|X}$.

$\langle \text{P:term-process} \rangle$ **Proposition 6.5.5.** *Suppose $\langle t_i \mid i \in I \rangle$ are jointly Borel given X . Then there exists a real-valued stochastic process, $\langle Y_i \mid i \in I \rangle$, defined on a probability space, (S, Γ, ν) , such that*

$$P(\bigwedge_{k=1}^n t_{i(k)} \in \underline{V}_k \mid X) = \nu \bigcap_{k=1}^n \{Y_{i(k)} \in V_k\}, \quad (6.5.1) \quad \boxed{\text{term-process}}$$

whenever $i(1), \dots, i(n) \in I$ and $V_k \in \mathcal{B}(\mathbb{R})$. Moreover, if $\langle t_i \mid i \in I \rangle$ are independent given X , then $\langle Y_i \mid i \in I \rangle$ are independent.

Proof. Let $S = \mathbb{R}^I$ and $\Gamma = \bigotimes_{i \in I} \mathcal{B}(\mathbb{R})$ the product σ -algebra. For each $i \in I$, let μ_i be the distribution of t_i . For each $n \in \mathbb{N}$ and $\mathbf{i} = (i(1), \dots, i(n)) \in I^n$, let $\mathbf{t}(\mathbf{i}) = (t_{i(1)}, \dots, t_{i(n)})$, so that $\mu_{\mathbf{t}(\mathbf{i})|X}$ is a probability measure on $(\mathbb{R}^n, \mathcal{B}(\mathbb{R}^n))$.

Using the methods in the proof of Proposition 6.5.1, we have that the measures $\mu_{\mathbf{t}(i)|X}$ are consistent, in the sense of Kolmogorov. Hence, by Kolmogorov's extension theorem (see, for instance, [15, Theorem 2.2.2]), there exists a probability measure ν on (S, Γ) such that

$$\nu\{x \in S \mid (x_{i(1)}, \dots, x_{i(n)}) \in A\} = \mu_{\mathbf{t}(i)|X} A, \quad (6.5.2) \quad \boxed{\text{term-process-2}}$$

whenever $A \in \mathcal{B}(\mathbb{R}^n)$. Define $Y_i : S \rightarrow \mathbb{R}$ by $Y_i(x) = x_i$. Then $\langle Y_i \mid i \in I \rangle$ is a real-valued stochastic process, and (6.5.1) follows from (6.5.2). The final claim about independence follows from Corollary 4.5.12, (6.5.2), and Remark 6.5.4. \square

6.5.4 The law of large numbers for terms

$\langle \text{L:LLN} \rangle$ **Lemma 6.5.6.** *Suppose $\langle t_n \mid n \in \mathbb{N} \rangle$ are jointly Borel given X , and let $\langle Y_n \mid n \in \mathbb{N} \rangle$ be as in Proposition 6.5.5. Define the terms $s_n = (t_1 + \dots + t_n)/n$ and the random variables $Z_n = (Y_1 + \dots + Y_n)/n$. Then*

$$P(\bigwedge_{k=1}^n s_{i(k)} \in \underline{V}_k \mid X) = \nu \bigcap_{k=1}^n \{Z_{i(k)} \in V_k\},$$

whenever $i(1), \dots, i(n) \in \mathbb{N}$ and $V_k \in \mathcal{B}(\mathbb{R})$.

Proof. Let $i(1), \dots, i(n) \in \mathbb{N}$ and $V_k \in \mathcal{B}(\mathbb{R})$. For each $n \in \mathbb{N}$, define $f_n : \mathbb{R}^n \rightarrow \mathbb{R}$ by $f_n(x_1, \dots, x_n) = (x_1 + \dots + x_n)/n$. Let $m = \max\{i(1), \dots, i(n)\}$ and $V'_k = f_{i(k)}^{-1} V_k \times \mathbb{R}^{m-i(k)}$. Then $\{Z_{i(k)} \in V_k\} = \{(Y_1, \dots, Y_m) \in V'_k\}$, so that $\bigcap_{k=1}^n \{Z_{i(k)} \in V_k\} = \{(Y_1, \dots, Y_m) \in V'\}$, where $V' = \bigcap_{k=1}^n V'_k$. By the definition of ν , we have $\nu \bigcap_{k=1}^n \{Z_{i(k)} \in V_k\} = \mu V'$, where μ is the distribution of (t_1, \dots, t_m) given X . Hence, by Remark 6.5.3,

$$\nu \bigcap_{k=1}^n \{Z_{i(k)} \in V_k\} = P((t_1, \dots, t_m) \in \underline{V}' \mid X).$$

But $\text{ZFC} \vdash \bigwedge_{k=1}^n s_{i(k)} \in \underline{V}_k \leftrightarrow (t_1, \dots, t_m) \in \underline{V}'$, so the result follows from the rule of logical implication and Proposition 3.2.14. \square

$\langle \text{T:LLN} \rangle$ **Theorem 6.5.7 (Law of large numbers).** *Let P be a real inductive theory in ZFC, and let $X \in \text{ante } P$. Let $\langle t_n \mid n \in \mathbb{N} \rangle$ be a sequence of terms that are i.i.d given X . Assume t_1 is integrable and let $\mu = E[t_1 \mid X]$. Define the terms $s_n = (t_1 + \dots + t_n)/n$, and let*

$$s = \{(x, y) \in \underline{\mathbb{N}}_0 \times \underline{\mathbb{R}} \mid \bigvee_{n \in \mathbb{N}_0} x = \underline{n} \wedge y = s_n\}.$$

Then $P(\varphi_{\lim}(s, \underline{\mu}) \mid X) = 1$.

Proof. By (6.3.2), we have $\text{ZFC}_\infty \vdash s \in \underline{\mathbb{R}}^{\underline{\mathbb{N}}_0}$ and $s(\underline{n}) = s_n$. Since we also have $\text{ZFC} \vdash \underline{\mu} \in \underline{\mathbb{R}}$, it follows that

$$\text{ZFC} \vdash \varphi_{\lim}(s, \underline{\mu}) \leftrightarrow (\forall z \in \underline{(0, \infty)}) \bigvee_{m=0}^\infty \bigwedge_{n=m}^\infty |s_n - \underline{\mu}| < z.$$

In fact, if we define $\varepsilon_\ell = 1/\ell$ for $\ell \in \mathbb{N}$, then

$$\text{ZFC} \vdash \varphi_{\lim}(s, \underline{\mu}) \leftrightarrow \bigwedge_{\ell=1}^{\infty} \bigvee_{m=0}^{\infty} \bigwedge_{n=m}^{\infty} |s_n - \underline{\mu}| < \varepsilon_\ell.$$

Using the real frame of reference, we have $\text{ZFC} \vdash |s_n - \underline{\mu}| < \varepsilon_\ell \leftrightarrow s_n \in \underline{V}_\ell$, where $\underline{V}_\ell = (\mu - \varepsilon_\ell, \mu + \varepsilon_\ell)$. Hence, $\text{ZFC} \vdash \varphi_{\lim}(s, \underline{\mu}) \leftrightarrow \psi$, where

$$\psi = \bigwedge_{\ell=1}^{\infty} \bigvee_{m=0}^{\infty} \bigwedge_{k=0}^{\infty} \bigwedge_{n=m}^{m+k} s_n \in \underline{V}_\ell.$$

Hence, by the rule of logical implication and Proposition 3.2.14, it suffices to show $P(\psi \mid X) = 1$.

By the continuity rule,

$$P(\psi \mid X) = \lim_{\ell \rightarrow \infty} \lim_{m \rightarrow \infty} \lim_{k \rightarrow \infty} P(\bigwedge_{n=m}^{m+k} s_n \in \underline{V}_\ell \mid X).$$

Let Y_n and Z_n be as in Lemma 6.5.6. Then the Y_n are i.i.d. and integrable, and $E^\nu[Y_1] = \mu$. By Lemma 6.5.6,

$$\begin{aligned} P(\psi \mid X) &= \lim_{\ell \rightarrow \infty} \lim_{m \rightarrow \infty} \lim_{k \rightarrow \infty} \nu \bigcap_{n=m}^{m+k} \{Z_n \in \underline{V}_\ell\} \\ &= \nu \bigcap_{\ell=1}^{\infty} \bigcup_{m=0}^{\infty} \bigcap_{n=m}^{\infty} \{|Z_n - \mu| < \varepsilon_\ell\}. \end{aligned}$$

On the other hand, by the law of large numbers for random variables, $Z_n \rightarrow \mu$ ν -a.s. Hence, $P(\psi \mid X) = 1$. \square

6.5.5 The central limit theorem for terms

Let t be Borel given X and assume that t is integrable. Let $r = E[t \mid X]$. Then the term $(t - \underline{r})^2 = (t - \underline{r}) \cdot (t - \underline{r})$ is Borel given X . The *variance of t given X* is defined by $V(t \mid X) = E[(t - \underline{r})^2 \mid X]$. Note that $V(t \mid X) \in [0, \infty]$.

Theorem 6.5.8 (Central limit theorem). *Let P be a real inductive theory in ZFC, and let $X \in \text{ante } P$. Let $\langle t_n \mid n \in \mathbb{N} \rangle$ be a sequence of terms that are i.i.d. given X . Assume t_1 is integrable. Let $\mu = E[t_1 \mid X]$ and let s_n be as in Theorem 6.5.7. Let $\sigma = \sqrt{V(t_1 \mid X)}$ and assume $\sigma \in (0, \infty)$. Then*

$$\lim_{n \rightarrow \infty} P(\sqrt{n} \cdot (s_n - \underline{\mu}) \leq \underline{r} \mid X) = \frac{1}{\sqrt{2\pi\sigma^2}} \int_{-\infty}^{\underline{r}} e^{-(x-\mu)^2/2\sigma^2} dx,$$

for all $r \in \mathbb{R}$.

Proof. Let Y_n and Z_n be as in Lemma 6.5.6. As in the proof of Theorem 6.5.7, we have

$$P(\sqrt{n} \cdot (s_n - \underline{\mu}) \leq \underline{r} \mid X) = \nu\{\sqrt{n}(Z_n - \mu) \leq \underline{r}\}.$$

The result therefore follows from the central limit theorem for random variables, applied to the sequence $\langle Y_n \mid n \in \mathbb{N} \rangle$. \square

6.6 Probabilities of probabilities

$\langle S: \text{cond-exp} \rangle$ Anyone who has played a tabletop role-playing game is familiar with the variety of dice that are used. Not only are there the usual 6-sided dice, but the standard collection also includes dice with 4, 8, 10, 12, and 20 sides. Imagine taking a collection of such dice, choosing one of them at random, and rolling it. In that case, it makes perfectly good sense to ask the question, “What is the probability that the probability of rolling a 1 is greater than 0.1?”

Although this example illustrates that is sensible to talk about probabilities of probabilities, it is a rather trivial example. More serious examples arise in applications of inference, where probabilities are updated on the basis of a succession of observations.

A priori, these nested probabilities seem to be a problem for inductive logic. An inductive statement is a triple, (X, φ, p) . The first two components of the triple are part of the language \mathcal{L} , but the triple itself is not. This makes it impossible to construct a formal sentence in \mathcal{L} that contains an inductive statement.

Measure-theoretic probability has a similar problem. There, probabilities are neither events nor random variables, so they cannot appear inside a probability measure. The way this is dealt with in measure theory is through the notion of conditional expectation. A conditional expectation *is*, in fact, a random variable, so it *can* appear inside a probability measure.

In this section, we construct the analogue of this for inductive logic. We will construct conditional probability and expectation, where we condition on terms and the resulting object is itself a term. We use this to formulate, in Theorem 6.6.8, the law of total probability, which is also known as the tower rule, or the law of iterated expectation.

6.6.1 Conditioning on terms

A *probability kernel* from \mathbb{R}^n to \mathbb{R}^m is a function $\nu : \mathbb{R}^n \times \mathcal{B}(\mathbb{R}^m) \rightarrow [0, 1]$ such that $\mu(\vec{r}, \cdot)$ is a probability measure for each $\vec{r} \in \mathbb{R}^n$, and $\mu(\cdot, V)$ is a measurable function for each $V \in \mathcal{B}(\mathbb{R}^m)$. If $m = n$, then we call μ a *probability kernel on* \mathbb{R}^n . If we say that μ is a probability kernel, or just a kernel, then we mean that μ is a probability kernel on \mathbb{R} .

$\langle D: \text{shrink-nice} \rangle$ **Definition 6.6.1.** Let ν be a probability measure on $(\mathbb{R}, \mathcal{B}(\mathbb{R}))$. For each $n \in \mathbb{N}_0$, let $V_n \in \mathcal{B}(\mathbb{R})$, and let $r \in \mathbb{R}$. Then V_n *shrinks nicely to* r with respect to ν if there exist $c > 0$ and $\varepsilon_n > 0$ such that $\varepsilon_n \rightarrow 0$, $V_n \subseteq (r - \varepsilon_n, r + \varepsilon_n)$, and $\nu V_n > c\nu(r - \varepsilon_n, r + \varepsilon_n)$.

$\langle T: \text{cond-dist} \rangle$ **Theorem 6.6.2.** Let P be a real inductive theory in ZFC, and let $X \in \text{ante } P$. Let $s, t \in \mathcal{T}$ be jointly Borel. Then there exists a probability kernel μ such that

$$\lim_{n \rightarrow \infty} P(t \in \underline{V} \mid X, s \in \underline{V}_n) = \mu(r, V) \quad (6.6.1) \quad \boxed{\text{cond-dist}}$$

for $\mu_{s|X}$ -a.e. $r \in \mathbb{R}$, whenever V_n shrinks nicely to r with respect to $\mu_{s|X}$. If $\tilde{\mu}$ is another such kernel, then $\mu(r, \cdot) = \tilde{\mu}(r, \cdot)$ for $\mu_{s|X}$ -a.e. $r \in \mathbb{R}$.

Proof. By Proposition 6.5.5, we may define random variables Y_1 and Y_2 on a probability space (S, Γ, ν) such that

$$P(s \in \underline{V} \wedge t \in \underline{V}' \mid X) = \nu\{Y_1 \in V\} \cap \{Y_2 \in V'\}$$

for all $V, V' \in \mathcal{B}(\mathbb{R})$.

Let \mathbf{E}^ν denote expectation with respect to ν , so that $\mathbf{E}^\nu[Z] = \int_S Z d\nu$ whenever Z is a real-valued random variable defined on (S, Γ, ν) . Then $\mathbf{E}^\nu[1_{\{Y_2 \in V\}} \mid Y_1]$ denotes a version of the conditional expectation of $1_{\{Y_2 \in V\}}$ given Y_1 . That is, $Z = \mathbf{E}^\nu[1_{\{Y_2 \in V\}} \mid Y_1]$ is a $\sigma(Y_1)$ -measurable random variable such that $\mathbf{E}^\nu[1_{\{Y_2 \in V\}} 1_B] = \mathbf{E}^\nu[Z 1_B]$ for all $B \in \sigma(Y_1)$. Moreover, if Z' is any other random variable with this property, then $Z = Z'$, ν -a.s.

There exists a probability kernel μ such that $\mathbf{E}^\nu[1_{\{Y_2 \in V\}} \mid Y_1] = \mu(Y_1, V)$, ν -a.s., for every $V \in \mathcal{B}(\mathbb{R})$. This kernel is unique in the sense that if $\tilde{\mu}$ is another such kernel, then there exists $N \in \Gamma$ such that $\nu\{Y_1 \in N\} = 0$ and $\mu(r, V) = \tilde{\mu}(r, V)$ for all $(r, V) \in N^c \times \mathcal{B}(\mathbb{R})$. The kernel μ is called a regular conditional distribution for Y_2 given Y_1 . The existence and uniqueness of regular conditional distributions is shown, for instance, in [14, Theorem 5.3].

Let V_n shrink nicely to r with respect to $\mu_{s|X}$. Then $P(s \in \underline{V}_n \mid X) = \mu_{s|X} V_n > 0$. By Proposition 6.5.5 and the multiplication rule,

$$P(t \in \underline{V} \mid X, s \in \underline{V}_n) = \frac{\nu\{Y_2 \in V\} \cap \{Y_1 \in V_n\}}{\nu\{Y_1 \in V_n\}}.$$

Using properties of conditional expectation and the fact that $\mu_{s|X} = \nu\{Y_1 \in \cdot\}$, we have

$$\begin{aligned} \nu\{Y_2 \in V\} \cap \{Y_1 \in V_n\} &= \mathbf{E}^\nu[1_{\{Y_1 \in V_n\}} \mathbf{E}^\nu[1_{\{Y_2 \in V\}} \mid Y_1]] \\ &= \mathbf{E}^\nu[1_{\{Y_1 \in V_n\}} \mu(Y_1, V)] \\ &= \int_{V_n} \mu(x, V) \mu_{s|X}(dx). \end{aligned}$$

Hence, by the Lebesgue differentiation theorem (see, for instance, [1, Theorem 8.4.6]),

$$P(t \in \underline{V} \mid X, s \in \underline{V}_n) = \frac{1}{\mu_{s|X}(V_n)} \int_{V_n} \mu(x, V) \mu_{s|X}(dx) \rightarrow \mu(r, V)$$

for $\mu_{s|X}$ -a.e. $r \in \mathbb{R}$.

Finally, suppose $\tilde{\mu}$ is another probability kernel such that (6.6.1) holds. Define $N \subseteq \mathbb{R}$ by $r \in N$ if and only if there exists $\varepsilon > 0$ such that $\mu_{s|X}(r - \varepsilon, r + \varepsilon) = 0$. Then N is an open set, and $\mu_{s|X} K = 0$ for all compact $K \subseteq N$. Hence, $\mu_{s|X} N = 0$. Let $r \in N^c$ and $V \in \mathcal{B}(\mathbb{R})$. Choose any $\varepsilon_n > 0$ with $\varepsilon_n \rightarrow 0$ and define $V_n = (r - \varepsilon_n, r + \varepsilon_n)$. Then $\mu_{s|X} V_n > 0$, so V_n satisfies Definition 6.6.1 with $c > 1/2$. Thus, V_n shrinks nicely to r with respect to $\mu_{s|X}$. By (6.6.1), we have $\mu(r, V) = \tilde{\mu}(r, V)$. \square

(C:cond-dist) **Corollary 6.6.3.** *Let P be a real inductive theory in ZFC, and let $X \in \text{ante } P$. Let $s, t \in \mathcal{T}$ be jointly Borel, and let μ be a probability kernel satisfying (6.6.1). Let $r \in \mathbb{R}$. If $P(t \in \underline{V} \mid X, s = \underline{r})$ exists, then*

$$P(t \in \underline{V} \mid X, s = \underline{r}) = \mu(r, V).$$

Proof. Choose any $\varepsilon_n > 0$ with $\varepsilon_n \rightarrow 0$ and let $V_n = (r - \varepsilon_n, r + \varepsilon_n)$. Suppose $P(t \in \underline{V} \mid X, s = \underline{r})$ exists. By Lemma 3.2.10, we have $P(s = \underline{r} \mid X) > 0$. Hence, $\mu_{s|X}\{r\} > 0$, so that V_n shrinks nicely to r with respect to $\mu_{s|X}$. By the multiplication and continuity rules, it follows that

$$\lim_{n \rightarrow \infty} P(t \in \underline{V} \mid X, s \in \underline{V}_n) = P(t \in \underline{V} \mid X, s = \underline{r}).$$

The result therefore follows from (6.6.1). \square

Remark 6.6.4. Theorem 6.6.2 and Corollary 6.6.3 can be generalized to the case where $s_1, \dots, s_n, t_1, \dots, t_m$ are jointly Borel and μ is a probability kernel from \mathbb{R}^n to \mathbb{R}^m . In that case, let $\mathbf{s} = (s_1, \dots, s_n)$, and replace (6.6.1) by

$$\lim_{n \rightarrow \infty} P(\bigwedge_{\ell=1}^m t_\ell \in \underline{V}_\ell \mid X, \bigwedge_{k=1}^n s_k \in \underline{V}_{k,n}) = \mu(r_1, \dots, r_n, V_1 \times \dots \times V_m)$$

for $\mu_{\mathbf{s}|X}$ -a.e. $(r_1, \dots, r_n) \in \mathbb{R}^n$, whenever $V_{k,n}$ shrinks nicely to r_k with respect to $\mu_{s_k|X}$.

Any kernel μ satisfying (6.6.1) is called a *distribution of t given X and s* .

6.6.2 Versions of distributions

For fixed V , a function $h(r)$ is called a *version of $P(t \in \underline{V} \mid X, s = \underline{r})$* if $h(r) = \mu(r, V)$ for some μ satisfying (6.6.1). Note that any two versions are equal $\mu_{s|X}$ -a.e. We may sometimes write $P(t \in \underline{V} \mid X, s = \underline{r}) = h(r)$ to mean that $h(r)$ is a version of $P(t \in \underline{V} \mid X, s = \underline{r})$, but it should be remembered that if $P(s = \underline{r} \mid X) = 0$, then $P(t \in \underline{V} \mid X, s = \underline{r})$ does not have a uniquely determined value for a fixed value of r .

Suppose t is integrable given X . A function $h(r)$ is called a *version of $E[t \mid X, s = \underline{r}]$* if $h(r) = \int_{\mathbb{R}} x \mu(r, dx)$ for some μ satisfying (6.6.1). Note that any two versions are equal $\mu_{s|X}$ -a.e. We may sometimes write $E[t \mid X, s = \underline{r}] = h(r)$ to mean that $h(r)$ is a version of $E[t \mid X, s = \underline{r}]$, but we should remember that if $P(s = \underline{r} \mid X) = 0$, then $E[t \mid X, s = \underline{r}]$ does not have a uniquely determined value for a fixed value of r .

Let μ satisfy (6.6.1). By the proof of Theorem 6.6.2 and properties of regular conditional distributions, we have $\int_{\mathbb{R}} |x| \mu(r, dx) < \infty$ for all $r \in \mathbb{R}$. It follows that $r \mapsto \int_{\mathbb{R}} x \mu(r, dx)$ is a version of $E[t \mid X, s = \underline{r}]$ whenever μ satisfies (6.6.1).

We may sometimes treat $P(t \in \underline{V} \mid X, s = \underline{r})$ and $E[t \mid X, s = \underline{r}]$ as if they were functions, rather than expressions that possess versions. In such situations, the meaning must be understood according to context. For example, if we say

that $P(t > 0 \mid X, s = \underline{r}) = 1 - e^{-r}$, then we mean that $h(r) = 1 - e^{-r}$ is a version of $P(t > 0 \mid X, s = \underline{r})$. On the other hand, if we say that

$$P(t > 0 \mid X, s = \underline{r}) = E[t' \mid X, s = \underline{r}],$$

then we mean that $h(r)$ is a version of $P(t > 0 \mid X, s = \underline{r})$ if and only if $h(r)$ is a version of $E[t' \mid X, s = \underline{r}]$. Since different versions are equal $\mu_{s|X}$ -a.e., such a claim can be verified by checking it for a single version.

All of this could be made precise if we were to define $E[t \mid X, s = \underline{r}]$ as the equivalence class of $h(r)$, where $h(r)$ is a version, and two functions are equivalent if they are equal $\mu_{s|X}$ -a.e. To save ourselves from even more notation, we avoid this approach. Moreover, it is common in measure-theoretic probability to use the language of “versions” when talking about conditional expectation.

6.6.3 Indicator terms

We have seen how to define $P(t \in \underline{V} \mid X, s = \underline{r})$. Here, we see how to define $P(\psi \mid X, s = \underline{r})$ for more general sentences ψ . If $\psi \in \mathcal{L}^0$, then we define the term $\underline{1}_\psi$ by

$$y = \underline{1}_\psi \leftrightarrow \psi \wedge y = \underline{1} \vee \neg\psi \wedge y = \underline{0}.$$

If $P(\psi \mid X)$ exists, then $\underline{1}_\psi$ is Borel given X . In fact, in this case, $\underline{1}_\psi$ is integrable and $E[\underline{1}_\psi \mid X] = P(\psi \mid X)$.

If s is Borel given X , then $\underline{1}_\psi$ and s are jointly Borel if and only if $P(\psi \wedge s \in \underline{V} \mid X)$ exists for all $V \in \mathcal{B}(\mathbb{R})$. In this case, we define

$$P(\psi \mid X, s = \underline{r}) = P(\underline{1}_\psi \in \underline{\{1\}} \mid X, s = \underline{r}).$$

That is, we say $h(r)$ is a version of $P(\psi \mid X, s = \underline{r})$ if and only if $h(r)$ is a version of $P(\underline{1}_\psi \in \underline{\{1\}} \mid X, s = \underline{r})$. Since $\text{ZFC} \vdash \psi \leftrightarrow \underline{1}_\psi \in \underline{\{1\}}$, Theorem 6.6.2 shows that

$$P(\psi \mid X, s \in \underline{V}_n) \rightarrow P(\psi \mid X, s = \underline{r}),$$

for $\mu_{s|X}$ -a.e. $r \in \mathbb{R}$, whenever V_n shrinks nicely to r with respect to $\mu_{s|X}$.

Note that

$$P(t \in \underline{V} \mid X, s = \underline{r}) = E[\underline{1}_\psi \mid X, s = \underline{r}],$$

where $\psi = t \in \underline{V}$. That is, $h(r)$ is a version of $P(t \in \underline{V} \mid X, s = \underline{r})$ if and only if $h(r)$ is a version of $E[\underline{1}_\psi \mid X, s = \underline{r}]$. Hence, distributions given X and s can be characterized entirely in terms of expectations.

The following result shows how to connect distributions given X and s to the root of our inductive theory.

Proposition 6.6.5. *Let P be a real inductive theory in ZFC with root T_0 , and let $X \in \text{ante } P$. Write $X \equiv T + \psi$, where $T \in [T_0, T_P]$. Let s and t be jointly Borel given T_0 and assume $P(\psi \wedge s \in \underline{V}_1 \wedge t \in \underline{V}_2 \mid T_0)$ exists for all $V_1, V_2 \in \mathcal{B}(\mathbb{R})$. Also assume that t is integrable given T_0 . Then $E[t \mid X, s = \underline{r}] = h(r, 1)$, where $h(r, r') = E[t \mid T_0, s = \underline{r}, \underline{1}_\psi = \underline{r}']$.*

Proof. Let $h(r, r')$ be a version of $E[t \mid T_0, s = \underline{r}, \underline{1}_\psi = \underline{r}']$. Let $\mathbf{s} = (s, \underline{1}_\psi)$. Then $h(r, r') = \int_{\mathbb{R}} x \mu(r, r', dx)$ for some kernel μ from \mathbb{R}^2 to \mathbb{R} such that

$$\lim_{n \rightarrow \infty} P(t \in \underline{V} \mid T_0, s \in \underline{V}_n, \underline{1}_\psi \in \underline{V}'_n) = \mu(r, r', V),$$

for $\mu_{\mathbf{s}|T_0}$ -a.e. $(r, r') \in \mathbb{R}^2$, whenever V_n shrinks nicely to r and V'_n shrinks nicely to r' . It follows that

$$\begin{aligned} \lim_{n \rightarrow \infty} P(t \in \underline{V} \mid X, s \in \underline{V}_n) &= \lim_{n \rightarrow \infty} P(t \in \underline{V} \mid T_0, \psi, s \in \underline{V}_n) \\ &= \lim_{n \rightarrow \infty} P(t \in \underline{V} \mid T_0, s \in \underline{V}_n, \underline{1}_\psi \in \underline{\{1\}}) \\ &= \mu(r, 1, V), \end{aligned}$$

for $\mu_{\mathbf{s}|X}$ -a.e. $r \in \mathbb{R}$, whenever V_n shrinks nicely to r . Hence, $h(r, 1) = \int_{\mathbb{R}} x \mu(r, 1, dx)$ is a version of $E[t \mid X, s = \underline{r}]$. \square

6.6.4 Conditional expectation

(P:cond-exp-versions) **Proposition 6.6.7.** *Let P be a real inductive theory in ZFC with root T_0 , and let $X \in \text{ante } P$. Let $s, t \in \mathcal{T}$ be jointly Borel, and assume t is integrable. If $h(r)$ and $h'(r)$ are versions of $E[t \mid X, s = \underline{r}]$, then $P(\underline{h}(s) = \underline{h}'(s) \mid X) = 1$.*

Proof. Let $V = \{r \in \mathbb{R} \mid h(r) = h'(r)\} \in \mathcal{B}(\mathbb{R})$. We first show that

$$\text{ZFC} \vdash \underline{h}(s) = \underline{h}'(s) \leftrightarrow s \in \underline{V}. \quad (6.6.2) \quad \boxed{\text{cond-exp-versions}}$$

Let $\mathcal{P} = (\Omega, \Sigma, \mathbb{P})$ be a model and assume $\mathcal{P} \models \text{ZFC}$ and $\mathcal{P} \models \underline{h}(s) = \underline{h}'(s)$. By adopting the real frame of reference, we may assume \mathcal{P} satisfies all of the conditions in Theorems 6.4.1, 6.4.3, and 6.4.5. Choose $\Omega^* \in \Sigma$ such that $\mathbb{P}\Omega^* = 1$ and, for every $\omega \in \Omega^*$, we have that Theorem 6.4.1(v) holds for V , Theorem 6.4.5 holds for both h and h' , and $\omega \models \underline{h}(s) = \underline{h}'(s)$. Let $\omega \in \Omega^*$. Then $s^\omega \in {}^\omega\mathbb{R}$ and there exists $b \in {}^\omega\mathbb{R}$ such that $b = h(s^\omega)$ and $b = h'(s^\omega)$. Hence, $s^\omega \in V \cap {}^\omega\mathbb{R} = {}^\omega\underline{V}$, so that $\omega \models s \in \underline{V}$. Since this is true for every $\omega \in \Omega^*$, we have $\mathcal{P} \models s \in \underline{V}$. Since \mathcal{P} was arbitrary, this gives $\text{ZFC}, \underline{h}(s) = \underline{h}'(s) \vdash s \in \underline{V}$. A similar proof shows that $\text{ZFC}, s \in \underline{V} \vdash \underline{h}(s) = \underline{h}'(s)$, and this verifies (6.6.2).

By (6.6.2) and Proposition 3.2.14, we have

$$P(\underline{h}(s) = \underline{h}'(s) \mid X) = P(s \in \underline{V} \mid X) = \mu_{\mathbf{s}|X} V = \mu_{\mathbf{s}|X} \{r \in \mathbb{R} \mid h(r) = h'(r)\}.$$

But h and h' are both versions of $E[t \mid X, s = \underline{r}]$. Therefore, $h = h'$, $\mu_{\mathbf{s}|X}$ -a.e., which shows that $P(\underline{h}(s) = \underline{h}'(s) \mid X) = 1$. \square

If $h(r)$ is a version of $E[t \mid X, s = \underline{r}]$, then the term $\underline{h}(s)$ is called a *version of $E[t \mid X, s]$* . Similarly, if $h(r)$ is a version of $P(t \in \underline{V} \mid X, s = \underline{r})$, then the term $\underline{h}(s)$ is called a *version of $P(t \in \underline{V} \mid X, s)$* . We call $E[t \mid X, s]$ the *conditional expectation of t given X and s* .

We may sometimes treat $E[t \mid X, s]$ as if it were a term, rather than an expression that possesses versions. In such situations, the meaning must be

understood according to context. For example, if we say that $t' = E[t \mid X, s]$, then we mean that t' is a version of $E[t \mid X, s]$. On the other hand, if we say that $E[E[t \mid X, s] \mid X] = E[t \mid X]$, then what we mean is that $E[t' \mid X] = E[t \mid X]$ whenever t' is a version of $E[t \mid X, s]$. A similar convention applies to versions of $P(t \in \underline{V} \mid X, s)$.

A version of $E[t \mid X, s = \underline{r}]$ is a measurable function, whereas a version of $E[t \mid X, s]$ is a term. Similarly, a version of $P(t \in \underline{V} \mid X, s = \underline{r})$ is a measurable function, whereas a version of $P(t \in \underline{V} \mid X, s)$ is a term. Terms can appear in inductive statements. Hence, for example,

$$P(P(t = \underline{1} \mid X, s) > \underline{0.5} \mid X) = 2/3 \quad (6.6.3) \quad \boxed{\text{cond-exp-expl}}$$

is a perfectly meaningful thing to say. It says that $(X, \underline{h}(s) > \underline{0.5}, 2/3) \in P$ whenever $\underline{h}(r)$ is a version of $P(t = \underline{1} \mid X, s = \underline{r})$. By Propositions 6.6.7 and 3.2.14, if we verify (6.6.3) for a single version, then it is true for all versions.

6.6.5 The law of total probability

(T:total-prob) **Theorem 6.6.8 (Law of total probability).** *Let P be a real inductive theory in ZFC and let $X \in \text{ante } P$. Let $s, t \in \mathcal{T}$ be jointly Borel, and assume t is integrable. Then $E[t \mid X] = E[E[t \mid X, s] \mid X]$.*

Proof. Let $\underline{h}(r)$ be a version of $E[t \mid X, s = \underline{r}]$. Then $\underline{h}(r) = \int_{\mathbb{R}} x \mu(r, dx)$ for some kernel μ satisfying (6.6.1). Let Y_1 and Y_2 be the random variables constructed in the proof of Theorem 6.6.2, so that μ is a regular conditional distribution for Y_2 given Y_1 . By properties of regular conditional distributions, we have $\mathbf{E}^\nu[Y_2 \mid Y_1] = \int_{\mathbb{R}} x \mu(Y_1, dx) = \underline{h}(Y_1)$. Using properties of conditional expectations, it follows that $\mathbf{E}^\nu[\underline{h}(Y_1)] = \mathbf{E}^\nu[\mathbf{E}^\nu[Y_2 \mid Y_1]] = \mathbf{E}^\nu[Y_2] = E[t \mid X]$. It therefore suffices to show that $E[s' \mid X] = \mathbf{E}^\nu[\underline{h}(Y_1)]$, where $s' = \underline{h}(s)$.

Using the real frame of reference, it follows easily that $\text{ZFC} \vdash s' \in \underline{V} \leftrightarrow s \in \underline{h}^{-1}\underline{V}$. Hence, s' is Borel given X and $\mu_{s'|X} = \mu_{s|X} \circ \underline{h}^{-1}$, which implies

$$E[s' \mid X] = \int_{\mathbb{R}} x \mu_{s'|X}(dx) = \int_{\mathbb{R}} \underline{h}(x) \mu_{s|X}(dx).$$

But $\mu_{s|X} = \nu\{Y_1 \in \cdot\}$, so $E[s' \mid X] = \mathbf{E}^\nu[\underline{h}(Y_1)]$. □

Suppose $P(\psi \wedge s \in \underline{V} \mid X)$ exists for each $V \in \mathcal{B}(\mathbb{R})$. In this case, we define $P(\psi \mid X, s) = E[\underline{1}_\psi \mid X, s]$, which means that a term t is a version of $P(\psi \mid X, s)$ if and only if it is a version of $E[\underline{1}_\psi \mid X, s]$. Since we also have that $P(\psi \mid X) = E[\underline{1}_\psi \mid X]$, Theorem 6.6.8 gives us

$$P(\psi \mid X) = E[P(\psi \mid X, s) \mid X].$$

This special case of Theorem 6.6.8, especially when $\mu_{s|X}$ is discrete, is what is more commonly known as the law of total probability.

Chapter 7

Principle of Indifference

(Ch:PoI) As discussed in Section 1.5, the principle of indifference is the heuristic idea that if we are equally ignorant about two statements, then we ought to assign them the same probability. The principle dates back to Laplace and the birth of mathematical probability. Although it is intuitively self-evident, it has a history of producing apparent paradoxes. It has no rigorous formulation in measure-theoretic probability theory. Hence, using measure theory alone, we are helpless to distinguish between valid and invalid uses of the principle.

In Section 7.1, we give a precise formulation of the principle of indifference in the context of inductive logic. As we will see, it is a natural generalization of a basic principle of deductive logic. In Sections 7.2 and 7.3, we present several elementary examples of the principle of indifference in action. All of these examples are finite, in the sense that they involves models $\mathcal{P} = (\Omega, \Sigma, \mathbb{P})$, where Ω is finite.

In the last three sections of this chapter, we consider the principle of indifference in the context of real inductive theories. In Section 7.4, we show that exchangeability is a special case of indifference. Sections 7.5 and 7.6 present several concrete examples. Section 7.5 treats examples involving an interval on the real line, while Section 7.6 treats examples involving circles and disks in the plane. The final example of Section 7.6 is the famous example of Bertrand's paradox.

7.1 Formulating the principle

(S:PoI-defn) The principle of indifference is an inductive generalization of a fundamental principle of deductive logic. This deductive principle is one that we use all the time, especially when we make assumptions “without loss of generality.” This principle, which we call “deductive indifference,” is presented in Section 7.1.2. In order to state it, we first define what we call “signature permutations.”

7.1.1 Signature permutations

Let L be a logical signature with associated predicate language \mathcal{L} . A *signature permutation*, or *L -permutation*, is a bijection $\pi : L \rightarrow L$ such that s^π has the same type and arity as s , for all $s \in L$. Given a signature permutation, we extend it to $\pi : \mathcal{T} \rightarrow \mathcal{T}$ by $x^\pi = x$ and $(ft_1 \cdots t_n)^\pi = f^\pi t_1^\pi \cdots t_n^\pi$. We then extend it to $\pi : \mathcal{L} \rightarrow \mathcal{L}$ by

- (i) $(s = t)^\pi = (s^\pi = t^\pi)$,
- (ii) $(rt_1 \cdots t_n)^\pi = r^\pi t_1^\pi \cdots t_n^\pi$,
- (iii) $(\neg \varphi)^\pi = \neg \varphi^\pi$,
- (iv) $(\bigwedge \Phi)^\pi = \bigwedge_{\varphi \in \Phi} \varphi^\pi$, and
- (v) $(\forall x \varphi)^\pi = \forall x \varphi^\pi$.

By the unique concatenation and reconstruction properties in Sections 5.1.1 and 5.1.2, each of these extensions of π is a bijection. Moreover, by formula induction in \mathcal{L}_{fin} , we have that $\varphi \in \mathcal{L}_{\text{fin}}$ if and only if $\varphi^\pi \in \mathcal{L}_{\text{fin}}$. For $X \subseteq \mathcal{L}$, we write $X^\pi = \{\varphi^\pi \mid \varphi \in X\}$. If $X^\pi \equiv X$, then we say that X is *invariant under π* , or *π -invariant*.

We will sometimes denote the inverse permutation, π^{-1} , by $-\pi$. Also, when π affects only finitely many extralogical symbols, we may use the usual cycle notation for permutations. For example, $\pi = (c_1 \ c_2)(r_1 \ r_2 \ r_3)$ means that $c_1^\pi = c_2$, $c_2^\pi = c_1$, $r_1^\pi = r_2$, $r_2^\pi = r_3$, $r_3^\pi = r_1$, and $s^\pi = s$ for all other extralogical symbols.

By term and formula induction, $\psi \in \text{Sf } \varphi$ if and only if $\psi^\pi \in \text{Sf } \varphi^\pi$. Also, $\text{var } t = \text{var } t^\pi$, and $\text{rk } \varphi$, $\text{var } \varphi$, $\text{bnd } \varphi$, and $\text{free } \varphi$ are all unchanged by replacing φ with φ^π . In particular, t is a ground term if and only if t^π is a ground term, and φ is a sentence if and only if φ^π is a sentence. On the other hand, $\text{sym } \varphi^\pi = \pi(\text{sym } \varphi)$ and $\text{con } \varphi^\pi = \pi(\text{con } \varphi)$.

If $\sigma : \text{Var} \rightarrow \mathcal{T}$ is a substitution, then define the substitution $\sigma' : \text{Var} \rightarrow \mathcal{T}$ by $\sigma' = \pi \circ \sigma$, which is the substitution given by $x^{\sigma'} = x^{\sigma\pi}$ for all $x \in \text{Var}$. Note that this relation does not extend to all of \mathcal{T} . For instance $c^{\sigma'} = c$, but $c^{\sigma\pi} = c^\pi$.

Proposition 7.1.1. *For all $t \in \mathcal{T}$ and all $\varphi \in \mathcal{L}$, we have $t^{\sigma\pi} = t^{\pi\sigma'}$ and $\varphi^{\sigma\pi} = \varphi^{\pi\sigma'}$. In particular, $\varphi(t/x)^\pi = \varphi^\pi(t^\pi/x)$.*

Proof. We first prove $t^{\sigma\pi} = t^{\pi\sigma'}$ by term induction. Since $x^\pi = x$, it is true for $x \in \text{Var}$ by the definition of σ' . Since constants are unaffected by substitutions, we have $c^{\sigma\pi} = c^\pi = c^{\pi\sigma'}$. Suppose it is true for t_1, \dots, t_n and let f be an n -ary function symbol. Then

$$\begin{aligned} (ft_1 \cdots t_n)^{\sigma\pi} &= (ft_1^\sigma \cdots t_n^\sigma)^\pi = f^\pi t_1^{\sigma\pi} \cdots t_n^{\sigma\pi} \\ &= f^\pi t_1^{\pi\sigma'} \cdots t_n^{\pi\sigma'} = (f^\pi t_1^\pi \cdots t_n^\pi)^{\sigma'} = (ft_1 \cdots t_n)^{\pi\sigma'}, \end{aligned}$$

and it is true for $ft_1 \cdots t_n$. By term induction, it holds for all $t \in \mathcal{T}$.

We next prove $\varphi^{\sigma\pi} = \varphi^{\pi\sigma'}$ by formula induction. The proof that it holds for prime formulas and for formulas of the form $\varphi = \neg\psi$ and $\varphi = \bigwedge \Phi$ is similar to the above. Suppose $\varphi = \forall x\psi$. Then $\varphi^{\sigma\pi} = (\forall x\psi^\tau)^\pi = \forall x\psi^{\tau\pi}$, where $x^\tau = x$ and $y^\tau = y^\sigma$ for $y \neq x$. By the inductive hypothesis, this gives $\varphi^{\sigma\pi} = \forall x\psi^{\pi\tau'}$. Using the fact that the proposition holds for terms, we have $x^{\tau'} = x$ and $y^{\tau'} = y^{\sigma'}$ for $y \neq x$. Hence, $\varphi^{\sigma\pi} = (\forall x\psi^\pi)^{\sigma'} = \varphi^{\pi\sigma'}$. The final assertion follows from the fact that if $\sigma = t/x$, then $\sigma' = t^\pi/x$. \square

Proposition 7.1.2. *Let $\varphi \in \mathcal{L}$. Then σ is free for φ if and only if σ' is free for φ^π . In particular, t is free for x in φ if and only if t^π is free for x in φ^π .*

Proof. First note that ζ is in the scope of $\forall z$ in φ if and only if ζ^π is in the scope of $\forall z$ in φ^π . Now suppose y is not free for x in φ . Then there exists $\zeta \in \text{Sf } \varphi$ such that $x \in \text{free } \zeta$, ζ is not in the scope of $\forall x$ in φ , and ζ is in the scope of $\forall y$ in φ . Since $\zeta^\pi \in \text{Sf } \varphi^\pi$ and $\text{free } \zeta^\pi = \text{free } \zeta$, it follows that y is not free for x in φ^π . The converse also holds. Hence, the second part of the proposition is true for $t = y$. For general t , simply note that $y \in t$ if and only if $y \in t^\pi$. The case of general σ now follows since σ is free for φ if and only if x^σ is free for x in φ for all x , and $x^{\sigma\pi} = x^{\sigma'}$. \square

7.1.2 Deductive indifference

(S:ded-PoI) Theorem 7.1.3 below says that the deductive derivability relation is preserved by signature permutations. The proof is elementary and the result is completely unsurprising. After all, the symbols that appear in a proof have no direct relevance. It is only their relationships to one another that matters. This principle is what we might call “deductive indifference.” As an example of applying this principle, we use it to formalize the technique of assuming something “without loss of generality.”

(T:pf-invar) **Theorem 7.1.3.** *Let $X \subseteq \mathcal{L}$ and $\varphi \in \mathcal{L}$. Let π be a signature permutation. Then $X^\pi \vdash \varphi^\pi$ if and only if $X \vdash \varphi$. In particular, if X is invariant under π , then $X \vdash \varphi^\pi$ if and only if $X \vdash \varphi$.*

Proof. Suppose $X \vdash \varphi$ and let $\langle \varphi_\beta \mid \beta \leq \alpha \rangle$ be a proof of φ from X . We claim that $\langle \varphi_\beta^\pi \mid \beta \leq \alpha \rangle$ is a proof of φ^π from X^π . Since $(\psi \rightarrow \zeta)^\pi = \psi^\pi \rightarrow \zeta^\pi$ and $(\bigwedge \Phi)^\pi = \bigwedge \Phi^\pi$, the claim will follow once we show that $\Lambda^\pi = \Lambda$.

For this, it suffices to show that $\Lambda \subseteq \Lambda^\pi$, since we can replace π by $-\pi$ and apply π to both sides. Showing $\Lambda \subseteq \Lambda^\pi$ can be done by verifying (I)–(IV) in Section 5.2.7, with Λ replaced by Λ^π . Verifying (II)–(IV) is straightforward. To show that $\Lambda^- \subseteq \Lambda^\pi$, it is enough to show $(\Lambda^-)^\pi \subseteq \Lambda$, for the reasons given above.

Let $\varphi \in \Lambda^-$. If φ has the form $(\Lambda 1)$, then so does φ^π , so that $\varphi^\pi \in \Lambda$. The same is true for $(\Lambda 2)$, $(\Lambda 3)$, $(\Lambda 5)$, and $(\Lambda 6)$. For $(\Lambda 4)$, suppose $\varphi = \forall x\psi \rightarrow \psi(t/x)$, where t is free for x in ψ . Then t^π is free for x in ψ^π , and $\varphi^\pi = \forall x\psi^\pi \rightarrow \psi^\pi(t^\pi/x)$. Hence, $\varphi^\pi \in \Lambda$. The proof for $(\Lambda 7)$ is similar.

This shows that $X \vdash \varphi$ implies $X^\pi \vdash \varphi^\pi$. Using the result with $-\pi$ gives the converse. The second result follows from the first by the definition of invariance. \square

Theorem 7.1.3 shows that if $T \subseteq \mathcal{L}^0$ is a deductive theory, then T^π is also. This is because if $T^\pi \vdash \varphi$, then $T \vdash \varphi^{-\pi}$, so that $\varphi^{-\pi} \in T$, which implies $\varphi \in T^\pi$. In fact, we have $T(X^\pi) = T(X)^\pi$ for any $X \subseteq \mathcal{L}^0$.

The following corollary is a formalization of the without-loss-of-generality proof method. It is also true in \mathcal{L}_{fin} if we require Φ to be finite.

$\langle \text{C:pf-invar} \rangle$ **Corollary 7.1.4.** *Let $X \subseteq \mathcal{L}$ and $\varphi \in \mathcal{L}$. Suppose $\Phi \subseteq \mathcal{L}$ is countable and $X \vdash \bigvee \Phi$. Fix $\theta_0 \in \Phi$ and suppose that for each $\theta \in \Phi$, there is a signature permutation π such that $\theta_0^\pi = \theta$, X is π -invariant, and φ is π -invariant. Then $X, \theta_0 \vdash \varphi$ implies $X \vdash \varphi$.*

Proof. Assume $X, \theta_0 \vdash \varphi$. Then $X \vdash \theta_0 \rightarrow \varphi$. But $\theta_0 \rightarrow \varphi \equiv \neg\varphi \rightarrow \neg\theta_0$. Hence, $X, \neg\varphi \vdash \neg\theta_0$. Let $\theta \in \Phi$. Choose π as in our hypotheses. Since $X, \neg\varphi$ is invariant under π , it follows from Theorem 7.1.3 that $X, \neg\varphi \vdash \neg\theta_0^\pi \equiv \neg\theta$. Since θ was arbitrary, $X, \neg\varphi \vdash \bigwedge_{\theta \in \Phi} \neg\theta \equiv \neg\bigvee \Phi$. Therefore, $X, \bigvee \Phi \vdash \varphi$. But $X \vdash \bigvee \Phi$, so we have $X \vdash \varphi$. \square

Example 7.1.5. Consider a scenario where we have three objects, and each is painted either red or blue. Define the extralogical signature $L = \{r, b\}$, where r and b are unary relation symbols. Let

$$X = \{\exists_{=3}, \forall x((rx \vee bx) \wedge \neg(rx \wedge bx))\}.$$

The set X asserts that there are three objects, each is red or blue, and none are both red and blue. Let us introduce the defined binary relation symbol s by

$$sxy \leftrightarrow (rx \wedge ry) \vee (bx \wedge by),$$

so that sxy asserts that x and y have the same color. We want to prove that there must be two objects of the same color. That is, we want to show that $X \vdash \varphi$, where

$$\varphi = \exists xy(x \neq y \wedge sxy).$$

Let us add a new constant c that represents an arbitrary object. By Proposition 5.2.16, it suffices to show that $X \vdash_{\mathcal{L}_c} \varphi$. Informally, we would like to say that, without loss of generality, we may assume c is red. In other words, we claim that it suffices to show $X, rc \vdash_{\mathcal{L}_c} \varphi$.

This is justified by Corollary 7.1.4. To see this, define π by $r^\pi = b$, $b^\pi = r$, and $c^\pi = c$. Then both X and φ are π -invariant. Let $\Phi = \{rc, bc\}$. Then $X \vdash rc \vee bc$ and $(rc)^\pi = bc$. Hence, according to Corollary 7.1.4, if we can establish $X, rc \vdash_{\mathcal{L}_c} \varphi$, then we may conclude $X \vdash_{\mathcal{L}_c} \varphi$.

7.1.3 Inductive indifference

We are now able to state the principle of indifference, which is an extension of Theorem 7.1.3 to the inductive setting.

Definition 7.1.6 (The Principle of Indifference). Let P be an inductive theory. Suppose that, for every signature permutation π , we have:

(R10) If $P(\varphi \mid X)$ exists and $X^\pi \in \text{ante } P$, then $P(\varphi^\pi \mid X^\pi) = P(\varphi \mid X)$.

Then P satisfies the *principle of indifference*.

If $P(\varphi \mid X)$ is to be an evidentiary relationship between X and φ , then the principle of indifference is a natural consistency condition to impose. After all, there is nothing special about the symbols we choose to use. In arithmetic, if we everywhere switch the symbols $+$ and \cdot , it is still the same arithmetic. The symbols just have the opposite meaning. If we were going to assign a certain probability in the original setting, then we ought to assign that same probability after the symbols are reversed, because nothing has actually changed.

In fact, it is such a natural requirement, it would have made sense to define it as one of our rules of inductive inference, making it part of the definition of an inductive theory. We did not do this for two reasons. First, the principle of indifference is a rule solely for predicate logic. Making it a rule of inductive inference would have created an asymmetry between propositional and predicate logic. Second, by omitting it from the rules of inductive inference, we were able to prove Theorem 5.4.2, which shows that all of modern, measure-theoretic probability is embedded in inductive logic. If we made the principle of indifference part of the definition of an inductive theory, this would not be the case. In other words, modern probability as we know it today, for better or for worse, does not require us to conform to the principle of indifference.

And so, the principle of indifference is not a required part of inductive logic. That is, inductive theories are required to satisfy (R1)–(R9), but they are not required to satisfy (R10). We will show, however, in the remaining sections of this chapter, how to add this requirement to our inferences by using inductive conditions.

Now, if X is π -invariant, then according to the rule of logical equivalence, we can reformulate the principle of indifference as $P(\varphi^\pi \mid X) = P(\varphi \mid X)$. This reformulation is perhaps closer to the intuitive idea of the principle of indifference. To say that X is our antecedent is to say that the totality of facts which we know consists of X , together with everything that can be proven from X . In other words, $T(X)$, the deductive theory generated by X , is the set of sentences that represents our knowledge. If X is π -invariant, then $T(X) = T(X^\pi)$. In other words, our knowledge remains entirely unchanged by the permutation π . In this sense, then, we cannot even “see” the permutation π . We are therefore “indifferent” between φ and φ^π . Everything we know about φ , we also know about φ^π , and vice versa. In this case, according to the principle of indifference, we should assign them the same probability.

7.1.4 Structures, models, and indifference

(S:PoI-models) In the remainder of this chapter, we will present several examples of the principle of indifference in action. In order to verify the principle of indifference in these examples, we will need several results about how (R10) relates to structures and models.

Let $\mathcal{P} = (\Omega, \Sigma, \mathbb{P})$ be an \mathcal{L} -model, and let π be a signature permutation. For $\omega = (A, L^\omega) \in \Omega$, define $\omega^\pi = (A, L^{\omega^\pi})$ so that $(s^\pi)^{\omega^\pi} = s^\omega$ for all $s \in L$. That is, $\omega^\pi = \omega \circ \pi^{-1}$. Let $\Omega^\pi = \{\omega^\pi \mid \omega \in \Omega\}$ and let $h_\pi : \Omega \rightarrow \Omega^\pi$ denote the map $\omega \mapsto \omega^\pi$. Let $\mathcal{P}^\pi = (\Omega^\pi, \Sigma^\pi, \mathbb{Q})$ be the measure space image of \mathcal{P} under h_π . Note that h_π is a bijection, and is therefore a pointwise isomorphism (as measure spaces) from \mathcal{P} to \mathcal{P}^π . Hence, it induces a measure-space isomorphism from $(\Omega, \Sigma, \mathbb{P})$ to $(\Omega^\pi, \Sigma^\pi, \mathbb{Q})$. In particular, $\mathbb{Q} = \mathbb{P} \circ h_\pi^{-1}$.

Given an assignment \mathbf{v} into \mathcal{P} , define the assignment \mathbf{v}^π into \mathcal{P}^π by $v_{\omega^\pi}^\pi(x) = v_\omega(x)$ for all $x \in \mathbf{Var}$. By term induction, we have $v_{\omega^\pi}^\pi(t^\pi) = v_\omega(t)$ for all $t \in \mathcal{T}$, and by formula induction, $\omega \models \varphi[v_\omega]$ if and only if $\omega^\pi \models \varphi^\pi[v_{\omega^\pi}^\pi]$, for any $\varphi \in \mathcal{L}$. Hence, $\varphi[\mathbf{v}]_\Omega = h_\pi^{-1} \varphi^\pi[\mathbf{v}^\pi]_{\Omega^\pi}$. Since $\mathbb{Q} = \mathbb{P} \circ h_\pi^{-1}$, this gives $\mathbb{P} \varphi[\mathbf{v}]_\Omega = \mathbb{Q} \varphi^\pi[\mathbf{v}^\pi]_{\Omega^\pi}$. In particular,

$$\mathcal{P} \models \varphi[\mathbf{v}] \text{ if and only if } \mathcal{P}^\pi \models \varphi^\pi[\mathbf{v}^\pi], \quad (7.1.1) \quad \boxed{\text{sem-pf-invar}}$$

for all $\varphi \in \mathcal{L}$ and all assignments \mathbf{v} into \mathcal{P} . Note that (7.1.1) could be used to give a semantic proof of Theorem 7.1.3.

(T:induc-invar) **Theorem 7.1.7.** *Let \mathcal{P} be an \mathcal{L} -model. Then for any $(X, \varphi, p) \in \mathcal{L}^{\text{IS}}$, we have $\mathcal{P} \models (X, \varphi, p)$ if and only if $\mathcal{P}^\pi \models (X^\pi, \varphi^\pi, p)$.*

Proof. Let $\mathcal{P} = (\Omega, \Sigma, \mathbb{P})$ be an \mathcal{L} -model and π a signature permutation, so that $\mathcal{P}^\pi = (\Omega^\pi, \Sigma^\pi, \mathbb{Q})$. Let $(X, \varphi, p) \in \mathcal{L}^{\text{IS}}$. Then $X \subseteq \mathcal{L}^0$ and $\varphi \in \mathcal{L}^0$. Suppose $\mathcal{P} \models (X, \varphi, p)$. Then there exists $Y \subseteq \mathcal{L}^0$ and $\psi \in \mathcal{L}^0$ such that $\mathcal{P} \models Y$, $X \equiv Y \cup \{\psi\}$, and $\mathbb{P} \varphi_\Omega \cap \psi_\Omega / \mathbb{P} \psi_\Omega = p$. By (7.1.1), we have $\mathcal{P}^\pi \models Y^\pi$. Theorem 7.1.3 implies $X^\pi \equiv Y^\pi \cup \{\psi^\pi\}$. And it follows from $\mathbb{Q} = \mathbb{P} \circ h_\pi^{-1}$ that $\mathbb{Q} \varphi_{\Omega^\pi}^\pi \cap \psi_{\Omega^\pi}^\pi / \mathbb{Q} \psi_{\Omega^\pi}^\pi = p$. Hence, $\mathcal{P}^\pi \models (X^\pi, \varphi^\pi, p)$. Applying this with $-\pi$ gives the converse. \square

Proposition 7.1.8. *Let \mathcal{P} be an \mathcal{L} -model. Suppose that $\mathcal{P} \simeq \mathcal{P}^\pi$ for every signature permutation π . Then $\mathbf{Th} \mathcal{P}$ satisfies the principle of indifference.*

Proof. Let $P = \mathbf{Th} \mathcal{P}$ and suppose $P(\varphi \mid X) = p$. Then $\mathcal{P} \models (X, \varphi, p)$, so that Theorem 7.1.7 gives $\mathcal{P}^\pi \models (X^\pi, \varphi^\pi, p)$. By hypothesis, $\mathcal{P} \simeq \mathcal{P}^\pi$. Hence, Theorem 5.3.24 implies $\mathcal{P} \models (X^\pi, \varphi^\pi, p)$. Therefore, $P(\varphi^\pi \mid X^\pi) = p$. \square

Proposition 7.1.9. *Let P be an inductive theory with root T_0 , and let $T \in [T_0, T_P]$. If P satisfies the principle of indifference, then so does $P \downarrow_{[T, T_P]}$.*

Proof. Let $P' = P \downarrow_{[T, T_P]}$. Proposition 3.5.10 implies that P' is an inductive theory. Suppose $P'(\varphi \mid X) = p$ and $X^\pi \in \text{ante } P'$. Since $P' \subseteq P$, we have $P(\varphi \mid X) = p$ and $X^\pi \in \text{ante } P$. By the principle of indifference for P , it follows that $P(\varphi^\pi \mid X^\pi) = p$. But $X^\pi \in \text{ante } P'$, and so we have $X^\pi \leftrightarrow [T, T_P]$. Therefore, $P'(\varphi^\pi \mid X^\pi) = p$. \square

$\langle \text{L:mapping-T}_0 \rangle$ **Lemma 7.1.10.** *Let P be an inductive theory with root T_0 and π a signature permutation. Let $X \subseteq \mathcal{L}^0$ and assume $X, X^\pi \in \text{ante } P$. Write $X \equiv T + \psi$ and $X^\pi \equiv T' + \psi'$, where $T, T' \in [T_0, T_P]$ and $\psi, \psi' \in \mathcal{L}^0$. Suppose $\mathcal{P} \models P$. Then, for all $\zeta \in T_0$, we have $\psi_\Omega \subseteq \zeta_\Omega^{-\pi}$ and $\psi'_\Omega \subseteq \zeta_\Omega^\pi$, \mathbb{P} -a.s. In particular, if $\zeta_\Omega^\pi \in \bar{\Sigma}$, then $\bar{\mathbb{P}} \zeta_\Omega^\pi > 0$, and if $\zeta_\Omega^{-\pi} \in \bar{\Sigma}$, then $\bar{\mathbb{P}} \zeta_\Omega^{-\pi} > 0$.*

Proof. Since $(\zeta^{-\pi})^\pi = \zeta \in T_0 \subseteq T(X^\pi) = T(X)^\pi$, we have $\zeta^{-\pi} \in T(X) \subseteq T_P + \psi$. Hence, $\psi \rightarrow \zeta^{-\pi} \in T_P$, which implies $\bar{\mathbb{P}} \psi_\Omega \cap (\zeta^{-\pi})_\Omega^c = 0$. Therefore, $\psi_\Omega \subseteq \zeta_\Omega^{-\pi}$, \mathbb{P} -a.s. In particular, if $\zeta_\Omega^{-\pi} \in \bar{\Sigma}$, then since $\bar{\mathbb{P}} \psi_\Omega > 0$, this gives $\bar{\mathbb{P}} \zeta_\Omega^{-\pi} > 0$. Similarly, since $\zeta \in T_0$, we have $\zeta^\pi \in T(X)^\pi = T(X^\pi) \subseteq T_P + \psi'$. Hence, $\psi' \rightarrow \zeta^\pi \in T_P$, so that $\psi'_\Omega \subseteq \zeta_\Omega^\pi$, \mathbb{P} -a.s. and $\zeta_\Omega^\pi \in \bar{\Sigma}$ implies $\bar{\mathbb{P}} \zeta_\Omega^\pi > 0$. \square

Let $\mathcal{P} = (\Omega, \Sigma, \mathbb{P})$ be a complete \mathcal{L} -model. Let $T_0 \subseteq \text{Th } \mathcal{P}$ and define $P = \mathbf{Th } \mathcal{P} \upharpoonright_{[T_0, \text{Th } \mathcal{P}]}$. Let π be a signature permutation. Suppose $P(\varphi \mid X) = p$ and $X^\pi \in \text{ante } P$. Let $\zeta \in T_0$ and assume $B = \zeta_\Omega^\pi \in \Sigma$. By Lemma 7.1.10, we have $\mathbb{P} B > 0$. We may therefore define the probability measures \mathbb{P}_B on (Ω, Σ) by $\mathbb{P}_B C = \mathbb{P} C \cap B / \mathbb{P} B$. Let $\mathcal{P}_B = (\Omega, \Sigma, \mathbb{P}_B)$. Note that $\mathcal{P}_B \models \text{Th } \mathcal{P}$. Similarly define $\mathcal{P}_{B'}$, where we assume $B' = \zeta_\Omega^{-\pi} \in \bar{\Sigma}$.

$\langle \text{P:mapping-T}_0 \rangle$ **Proposition 7.1.11.** *With the notation given above, if $\mathcal{P}_B \simeq \mathcal{P}_{B'}^\pi$, then $P(\varphi^\pi \mid X^\pi) = p$.*

Proof. Assume $\mathcal{P}_B \simeq \mathcal{P}_{B'}^\pi$. Write $X \equiv T + \psi$ and $X^\pi \equiv T' + \psi'$, where $T, T' \in [T_0, \text{Th } P]$ and $\psi, \psi' \in \mathcal{L}^0$. Since $P(\varphi \mid X) = p$, we have $\mathbb{P} \varphi_\Omega \cap \psi_\Omega / \mathbb{P} \psi_\Omega = p$. But $\psi_\Omega \subseteq B'$ a.s., by Lemma 7.1.10, so that

$$p = \frac{\mathbb{P} \varphi_\Omega \cap \psi_\Omega \cap B'}{\mathbb{P} \psi_\Omega \cap B'} = \frac{\mathbb{P}_{B'} \varphi_\Omega \cap \psi_\Omega}{\mathbb{P}_{B'} \psi_\Omega},$$

which implies $\mathcal{P}_{B'} \models (X, \varphi, p)$. By Theorems 7.1.7 and 5.3.24, it follows that $\mathcal{P}_B \models (X^\pi, \varphi^\pi, p)$. As above, since $\psi'_\Omega \subseteq B$, this gives

$$p = \frac{\mathbb{P}_B \varphi_\Omega^\pi \cap \psi'_\Omega}{\mathbb{P}_B \psi'_\Omega} = \frac{\mathbb{P} \varphi_\Omega^\pi \cap \psi'_\Omega \cap B}{\mathbb{P} \psi'_\Omega \cap B} = \frac{\mathbb{P} \varphi_\Omega^\pi \cap \psi'_\Omega}{\mathbb{P} \psi'_\Omega}.$$

Therefore, $\mathcal{P} \models (X^\pi, \varphi^\pi, p)$, so that $P(\varphi^\pi \mid X^\pi) = p$. \square

7.2 Examples with a single object

$\langle \text{S:basic-expls-PoI} \rangle$ In this section, we give several elementary examples involving the principle of indifference. In each of these examples, we are primarily concerned with the properties of a single object.

7.2.1 Either it's true or it isn't

The most naive misapplication of the principle of indifference is illustrated by the following invalid reasoning. Imagine we have a book whose color is unknown

to us. It might be red or might not be red. Since we have no reason to think one way or the other, we should assign equal probabilities to both cases. Therefore, the probability the book is red is $1/2$.

Clearly, this cannot be correct, for we could also apply it to the color black, and then to the color blue. Since probabilities must add up to one, it cannot be the case that all three colors have probability $1/2$.

To see formally that this argument is invalid, let $L = \{b, R\}$, where b is a constant symbol that denotes the book, and R is a unary predicate symbol that denotes the property of being red. Recall the notation, \mathcal{J}_{T_0} , for the set of inductive theories with root T_0 . We will assume only the principle of indifference. That is, we take $T_0 = \mathbf{Taut}$, and define the inductive condition,

$$\mathcal{C} = \{P \in \mathcal{J}_{\mathbf{Taut}} \mid P(R(b) \mid \mathbf{Taut}) \text{ exists and } P \text{ satisfies (R10)}\}.$$

Proposition 7.2.1. *The condition \mathcal{C} is consistent and $\mathcal{C} \not\vdash (\mathbf{Taut}, R(b), 1/2)$.*

Proof. Let $A = \{0\}$ and define $\omega_0 = (A, L^{\omega_0})$ by $b^{\omega_0} = 0$ and $R^{\omega_0} = \emptyset$. Define ω_1 similarly, but with $R^{\omega_1} = \{0\}$. Let $\Omega = \{\omega_0, \omega_1\}$ and $\Sigma = \mathfrak{P}\Omega$. Fix $p \in (0, 1)$ and define $\mathcal{P} = (\Omega, \Sigma, \mathbb{P})$, where $\mathbb{P}\{\omega_0\} = 1 - p$ and $\mathbb{P}\{\omega_1\} = p$. Let $P = \mathbf{Th} \mathcal{P} \in \mathcal{J}_{\mathbf{Taut}}$. Note that $R(b)_\Omega = \{\omega_1\}$, so that $\mathbb{P} R(b)_\Omega = p$. Thus, $P(R(b) \mid \mathbf{Taut}) = p$.

We will show that P satisfies (R10). Suppose that $P(\varphi \mid X) = p$ and $X^\pi \in \text{ante } P$. In this example, the only signature permutation is the identity. Hence, $P(\varphi^\pi \mid X^\pi) = p$, and so P satisfies (R10). This shows that $P \in \mathcal{C}$ and therefore \mathcal{C} is consistent. However, since $P(R(b) \mid \mathbf{Taut}) = p$ and p was arbitrary, it follows that $(\mathbf{Taut}, R(b), 1/2) \notin \mathbf{P}_{\mathcal{C}}$. \square

7.2.2 A single coin flip

$\langle \mathbf{S} : \text{one-coin} \rangle$ Let us return to the example in Section 5.4.3. We flip a coin with two sides, and assume only that the sides are distinct and that the coin will land on one of them. We also assume the principle of indifference.

Let $L = \{c, s_0, s_1\}$, where c , s_0 , and s_1 are constant symbols. We think of s_1 and s_0 as denoting the heads and tails sides of the coin, respectively, and c as denoting the result of our toss. Let T_0 be generated by the sentences

$$\begin{aligned} \varphi_1 &: s_0 \neq s_1 \\ \varphi_2 &: c = s_0 \vee c = s_1 \end{aligned}$$

Let $\mathcal{C} = \{P \in \mathcal{J}_{T_0} \mid P(c = s_1 \mid T_0) \text{ exists and } P \text{ satisfies (R10)}\}$.

$\langle \mathbf{P} : \text{one-coin} \rangle$ **Proposition 7.2.2.** *The condition \mathcal{C} is consistent and $\mathbf{P}_{\mathcal{C}}(c = s_1 \mid T_0) = 1/2$.*

Proof. Let $A = \{0, 1\}$ and define $\omega_0 = (A, L^{\omega_0})$ by $s_i^{\omega_0} = i$ and $c^{\omega_0} = 0$. Define ω_1 similarly, but with $c^{\omega_1} = 1$. Let $\Omega = \{\omega_0, \omega_1\}$, $\Sigma = \mathfrak{P}\Omega$, and define $\mathcal{P} = (\Omega, \Sigma, \mathbb{P})$, where $\mathbb{P}\{\omega_0\} = \mathbb{P}\{\omega_1\} = 1/2$. Since $\mathcal{P} \models T_0$, we may define $P = \mathbf{Th} \mathcal{P} \upharpoonright_{[T_0, \text{Th} \mathcal{P}]} \in \mathcal{J}_{T_0}$. Note that $P(c = s_1 \mid T_0) = 1/2$.

We will show that P satisfies (R10). Suppose that $P(\varphi \mid X) = p$ and $X^\pi \in \text{ante } P$. First assume $c^\pi = \mathbf{s}_0$. Then $\mathbf{s}_i^\pi = c$ for some $i \in \{0, 1\}$. Let $i' = 1 - i$, so that $\mathbf{s}_{i'}^\pi = \mathbf{s}_1$. We will apply Proposition 7.1.11. Let $\zeta = (\mathbf{s}_0 \neq \mathbf{s}_1)$. Then

$$\begin{aligned} B &= \zeta_\Omega^\pi = (c \neq \mathbf{s}_1)_\Omega = \{\omega_0\}, \text{ and} \\ B' &= \zeta_\Omega^{-\pi} = (c \neq \mathbf{s}_{i'})_\Omega = \{\omega_i\}. \end{aligned}$$

Define $h : \mathcal{P}_B \rightarrow \mathcal{P}_{B'}^\pi$ by $\omega_0 \mapsto \omega_i^\pi$ and $\omega_1 \mapsto \omega_{i'}^\pi$. Since $\mathbb{P}_B\{\omega_0\} = 1$ and $\mathbb{P}_{B'}^\pi\{\omega_i^\pi\} = \mathbb{P}_{B'}\{\omega_i\} = 1$, the function h induces a measure space isomorphism.

Let $g : A \rightarrow A$ be the bijection defined by $g(0) = i$ and $g(1) = i'$. Note that if $i = 0$, then $\pi = (c \ \mathbf{s}_0)$, and if $i = 1$, then $\pi = (c \ \mathbf{s}_0 \ \mathbf{s}_1)$. In either case, $g \circ \omega_0 = \omega_i \circ \pi^{-1} = \omega_i^\pi$. Hence, ω_i^π is the isomorphic image of ω_0 under g . Since $\mathbb{P}_B\{\omega_0\} = 1$, we have $\omega \simeq h\omega$, \mathbb{P}_B -a.s., so that h is a model isomorphism. Proposition 7.1.11 therefore gives $P(\varphi^\pi \mid X^\pi) = p$. The case $c^\pi = \mathbf{s}_1$ is similar.

Now assume $c^\pi = c$. If π is the identity, then $P(\varphi \mid X) = P(\varphi^\pi \mid X^\pi)$. Assume $\pi = (\mathbf{s}_0 \ \mathbf{s}_1)$. Define $h : \mathcal{P} \rightarrow \mathcal{P}^\pi$ by $\omega_i \mapsto \omega_{i'}^\pi$. Since $\mathbb{P}\{\omega_i\} = 1/2$ and $\mathbb{P}^\pi\{\omega_{i'}^\pi\} = \mathbb{P}\{\omega_{i'}\} = 1/2$, the function h induces a measure space isomorphism. As above, if $g : A \rightarrow A$ is given by $g(i) = i'$, then $\omega_{i'}^\pi = \omega_{i'} \circ \pi^{-1} = g \circ \omega_i$. Therefore, $\omega \simeq h\omega$, \mathbb{P} -a.s., so that h is a model isomorphism. By Theorems 7.1.7 and 5.3.24, this gives $P(\varphi^\pi \mid X^\pi) = p$. Altogether, this shows P satisfies the principle of indifference, and so $P \in \mathcal{C}$. Therefore, \mathcal{C} is consistent.

For the final claim, it suffices to show that $P(c = \mathbf{s}_1 \mid T_0) = 1/2$ whenever $P \in \mathcal{C}$. Let $P \in \mathcal{C}$ be given. Then P is an inductive theory with root T_0 , $P(c = \mathbf{s}_1 \mid T_0) = p$ for some p , and P satisfies the principle of indifference. Let $\pi = (\mathbf{s}_0 \ \mathbf{s}_1)$, and note that $T_0^\pi = T_0$. Hence, by the principle of indifference, we have $P(c = \mathbf{s}_0 \mid T_0) = p$. But $T_0 \vdash \neg(c = \mathbf{s}_0 \wedge c = \mathbf{s}_1)$, so that the addition rule gives $P(c = \mathbf{s}_0 \vee c = \mathbf{s}_1 \mid T_0) = 2p$. On the other hand, $T_0 \vdash c = \mathbf{s}_0 \vee c = \mathbf{s}_1$, so that the rule of logical implication implies $P(c = \mathbf{s}_0 \vee c = \mathbf{s}_1 \mid T_0) = 1$. By the rule of logical equivalence, we must have $2p = 1$, or $p = 1/2$. Thus, $P(c = \mathbf{s}_1 \mid T_0) = 1/2$. \square

7.2.3 A single trial

(S:one-trial) Here we consider a single “trial” that can result in either success or failure. Intuitively, this example is the same as the single coin flip in Section 7.2.2. With the coin flip, the possible results (heads or tails) were represented by objects. In this example, the possible results (success or failure) will be represented by predicates. For this reason, the proof of consistency in this example will be shorter.

Let $L = \{\mathbf{t}, S, F\}$, where \mathbf{t} is a constant symbol and S and F are unary relation symbols. We think of \mathbf{t} as denoting the trial, and S and F as denoting the properties of success and failure, respectively.

Let T_0 be generated by the sentence

$$\varphi : \forall x((Sx \vee Fx) \wedge \neg(Sx \wedge Fx))$$

Let $\mathcal{C} = \{P \in \mathfrak{I}_{T_0} \mid P(S(\mathbf{t}) \mid T_0) \text{ exists and } P \text{ satisfies (R10)}\}$.

Proposition 7.2.3. *The condition \mathcal{C} is consistent and $\mathbf{P}_{\mathcal{C}}(S(\mathbf{t}) \mid T_0) = 1/2$.*

Proof. Let $A = \{0, 1\}$ and define $\omega_0 = (A, L^{\omega_0})$ by $S^{\omega_0} = \{1\}$, $F^{\omega_0} = \{0\}$, and $\mathbf{t}^{\omega_0} = 0$. Define ω_1 similarly, but with $\mathbf{t}^{\omega_1} = 1$. Let $\Omega = \{\omega_0, \omega_1\}$, $\Sigma = \mathfrak{P}\Omega$, and define $\mathcal{P} = (\Omega, \Sigma, \mathbb{P})$, where $\mathbb{P}\{\omega_0\} = \mathbb{P}\{\omega_1\} = 1/2$. Since $\mathcal{P} \models T_0$, we may define $P = \mathbf{Th} \mathcal{P} \downarrow_{[T_0, Th \mathcal{P}]} \in \mathfrak{I}_{T_0}$. Note that $P(S(\mathbf{t}) \mid T_0) = 1/2$.

We will show that P satisfies (R10). Suppose that $P(\varphi \mid X) = p$ and $X^\pi \in \text{ante } P$. If π is the identity, then $P(\varphi \mid X) = P(\varphi^\pi \mid X^\pi)$. If π is not the identity, then $\pi = (S F)$. Define $h : \mathcal{P} \rightarrow \mathcal{P}^\pi$ by $\omega_i \mapsto \omega_{i'}^\pi$, where $i' = 1 - i$. Since $\mathbb{P}\{\omega_i\} = 1/2$ and $\mathbb{P}^\pi\{\omega_{i'}^\pi\} = \mathbb{P}\{\omega_{i'}\} = 1/2$, the function h induces a measure space isomorphism. If $g : A \rightarrow A$ is given by $g(i) = i'$, then $\omega_{i'}^\pi = \omega_{i'} \circ \pi^{-1} = g \circ \omega_i$, so that $\omega_i \simeq \omega_{i'}^\pi$. Therefore, $\omega \simeq h\omega$, \mathbb{P} -a.s., and h is a model isomorphism. By Theorems 7.1.7 and 5.3.24, this gives $P(\varphi^\pi \mid X^\pi) = p$. Altogether, this shows P satisfies the principle of indifference, and so $P \in \mathcal{C}$. Therefore, \mathcal{C} is consistent. The proof that $\mathbf{P}_{\mathcal{C}}(S(\mathbf{t}) \mid T_0) = 1/2$ follows as in the proof of Proposition 7.2.2. \square

7.2.4 Success is good

(S:success-good) In this example, we again consider a single trial that can result in either success or failure. This time, however, we include the qualitative information that success is “good.” This produces an asymmetry between success and failure. We will see that, because of this asymmetry, we can no longer conclude that the probability of success is $1/2$.

Let $L = \{\mathbf{t}, S, F, G\}$, where \mathbf{t} is a constant symbol, and S, F , and G are unary relation symbol. We think of \mathbf{t} as denoting the trial, S and F as denoting the properties of success and failure, and G as denoting the property of “goodness.”

Let T_0 be generated by the sentences

$$\begin{aligned}\varphi_1 &: \forall x((Sx \vee Fx) \wedge \neg(Sx \wedge Fx)) \\ \varphi_2 &: \forall x(Sx \rightarrow Gx)\end{aligned}$$

Let $\mathcal{C} = \{P \in \mathfrak{I}_{T_0} \mid P(S(\mathbf{t}) \mid T_0) \text{ exists and } P \text{ satisfies (R10)}\}$.

Proposition 7.2.4. *The condition \mathcal{C} is consistent and $\mathcal{C} \not\vdash (T_0, S(\mathbf{t}), 1/2)$.*

Proof. Let $A = \{0, 1\}$ and define $\omega_0 = (A, L^{\omega_0})$ by $S^{\omega_0} = G^{\omega_0} = \{1\}$, $F^{\omega_0} = \{0\}$, and $\mathbf{t}^{\omega_0} = 0$. Define ω_1 similarly, but with $\mathbf{t}^{\omega_1} = 1$. Let $\Omega = \{\omega_0, \omega_1\}$ and $\Sigma = \mathfrak{P}\Omega$. Fix $p \in (0, 1)$ and define $\mathcal{P} = (\Omega, \Sigma, \mathbb{P})$, where $\mathbb{P}\{\omega_0\} = 1 - p$ and $\mathbb{P}\{\omega_1\} = p$. Since $\mathcal{P} \models T_0$, we may define $P = \mathbf{Th} \mathcal{P} \downarrow_{[T_0, Th \mathcal{P}]} \in \mathfrak{I}_{T_0}$. Note that $P(S(\mathbf{t}) \mid T_0) = p$.

We will show that P satisfies (R10). Suppose that $P(\varphi \mid X) = p$ and $X^\pi \in \text{ante } P$. We may assume that π is not the identity permutation. Since there is only one constant symbol, we have $\mathbf{t}^\pi = \mathbf{t}$. First assume $G^\pi = F$. Let $\zeta = \forall x(Sx \rightarrow Gx)$. Then $\zeta^\pi = \forall x(S^\pi x \rightarrow Fx)$. By Lemma 7.1.10, we must have $\zeta_\Omega^\pi \neq \emptyset$, and so $S^\pi \neq S$. Therefore, $S^\pi = G$ and $F^\pi = S$. But this implies

$\zeta^{-\pi} = \forall x(Fx \rightarrow Sx)$, so that $\overline{\mathbb{P}}\zeta_{\Omega}^{-\pi} = 0$, contradicting Lemma 7.1.10. Hence, $G^{\pi} \neq F$. By reversing the roles of π and π^{-1} in this argument, we may also conclude that $F^{\pi} \neq G$.

Now assume $S^{\pi} = F$, so that $F^{\pi} = S$ and $G^{\pi} = G$. Then $\zeta^{\pi} = \forall x(Fx \rightarrow Gx)$, which gives $\overline{\mathbb{P}}\zeta_{\Omega}^{\pi} = 0$, again contradicting Lemma 7.1.10. Therefore, $S^{\pi} \neq F$. It now follows that $F^{\pi} = F$, $S^{\pi} = G$, and $G^{\pi} = S$. But this implies $\omega_i^{\pi} = \omega_i$, so that $\mathcal{P}^{\pi} = \mathcal{P}$. Thus, by Theorem 7.1.7, we have $P(\varphi^{\pi} \mid X^{\pi}) = p$, and P satisfies the principle of indifference. This shows that $P \in \mathcal{C}$, so that \mathcal{C} is consistent. Since p was arbitrary, we have $\mathcal{C} \not\vdash (T_0, S(\mathbf{t}), 1/2)$. \square

7.2.5 Goodness is independent

In Section 7.2.3, we considered a single trial, which could result in success or failure. There, we used the principle of indifference to conclude that the probability of success was $1/2$. In Section 7.2.4, we added the assumption that success is good. By doing so, we were no longer able to use the principle of indifference.

This may feel counterintuitive. To assert that success is good seems quite natural, and we might not expect this to invalidate our use of the principle of indifference. One reason we might feel this way is that we may have an ingrained sense that the property of goodness should not affect success or failure. But if this is a fact we wish to assume, then we must make it explicit. In this example, we will do exactly that. We will assume that whether or not \mathbf{t} is a success is independent of the fact that success is a good outcome.

Let $L = \{\mathbf{t}, S, F, G\}$ as in Section 7.2.4. Let T_0 be generated by the sentence

$$\varphi : \forall x((Sx \vee Fx) \wedge \neg(Sx \wedge Fx))$$

Let \mathcal{C} be the set of $P \in \mathfrak{J}_{T_0}$ such that

- (i) $P(S(\mathbf{t}) \mid T_0)$ exists,
- (ii) P satisfies the principle of indifference, and
- (iii) $S(\mathbf{t})$ and $\forall x(Sx \rightarrow Gx)$ are independent given T_0 .

$\langle \text{P:one-trial-indep} \rangle$ **Proposition 7.2.5.** *The condition \mathcal{C} is consistent and*

$$\mathbf{P}_{\mathcal{C}}(S(\mathbf{t}) \mid T_0, \forall x(Sx \rightarrow Gx)) = 1/2. \quad (7.2.1) \quad \boxed{\text{one-trial-indep}}$$

Proof. Let $A = \{0, 1\}$ and define $\omega_0 = (A, L^{\omega_0})$ by $S^{\omega_0} = \{1\}$, $F^{\omega_0} = \{0\}$, $G^{\omega_0} = A$, and $\mathbf{t}^{\omega_0} = 0$. Define ω_1 similarly, but with $\mathbf{t}^{\omega_1} = 1$. Let $\Omega = \{\omega_0, \omega_1\}$, $\Sigma = \mathfrak{P}\Omega$, and define $\mathcal{P} = (\Omega, \Sigma, \mathbb{P})$, where $\mathbb{P}\{\omega_0\} = \mathbb{P}\{\omega_1\} = 1/2$. Since $\mathcal{P} \models T_0$, we may define $P = \mathbf{Th} \mathcal{P} \upharpoonright_{[T_0, Th \mathcal{P}]} \in \mathfrak{J}_{T_0}$. Note that $P(S(\mathbf{t}) \mid T_0) = 1/2$, so that (i) holds. Also, $P(\forall x(Sx \rightarrow Gx) \mid T_0) = 1$, so that (iii) holds.

We will show that P satisfies (R10). Suppose that $P(\varphi \mid X) = p$ and $X^{\pi} \in \text{ante } P$. We may assume that π is not the identity permutation. Since there is only one constant symbol, we have $\mathbf{t}^{\pi} = \mathbf{t}$. First assume $S^{\pi} = G$. Let

$\zeta = \forall x \neg(Sx \wedge Fx)$. Then $\zeta^\pi = \forall x \neg(Gx \wedge F^\pi x)$. Since, for all $\omega \in \Omega$, we have $G^\omega = A$ and $(F^\pi)^\omega \neq \emptyset$, this gives $\zeta_\Omega^\pi = \emptyset$, contradicting Lemma 7.1.10. Hence, $S^\pi \neq G$. Similarly, $F^\pi \neq G$. We must therefore have $\pi = (S F)$.

Define $h : \mathcal{P} \rightarrow \mathcal{P}^\pi$ by $\omega_i \mapsto \omega_{i'}^\pi$, where $i' = 1 - i$. Since $\mathbb{P}\{\omega_i\} = 1/2$ and $\mathbb{P}^\pi\{\omega_{i'}^\pi\} = \mathbb{P}\{\omega_{i'}\} = 1/2$, the function h induces a measure space isomorphism. If $g : A \rightarrow A$ is given by $g(i) = i'$, then $\omega_{i'}^\pi = \omega_{i'} \circ \pi^{-1} = g \circ \omega_i$, so that $\omega_i \simeq \omega_{i'}^\pi$. Therefore, $\omega \simeq h\omega$, \mathbb{P} -a.s., and h is a model isomorphism. By Theorems 7.1.7 and 5.3.24, this gives $P(\varphi^\pi \mid X^\pi) = p$. Altogether, this shows P satisfies the principle of indifference, and so $P \in \mathcal{C}$. Therefore, \mathcal{C} is consistent.

Now let $P \in \mathcal{C}$ be arbitrary. Let $\pi = (S F)$. Then $T_0^\pi = T_0$, so that $P(S(\mathbf{t}) \mid T_0) = P(F(\mathbf{t}) \mid T_0)$, which implies $P(S(\mathbf{t}) \mid T_0) = 1/2$. Since P was arbitrary, we have $\mathbf{P}_\mathcal{C}(S(\mathbf{t}) \mid T_0) = 1/2$. Therefore, (iii) and the definition of independence give (7.2.1). \square

7.2.6 Lowering the root

We will take one last look at the example of the single trial. As before, let $L = \{\mathbf{t}, S, F, G\}$ as in Section 7.2.4. Let T_0 be generated by the sentence

$$\varphi : \forall x((Sx \vee Fx) \wedge \neg(Sx \wedge Fx))$$

Let \mathcal{C} be the set of $P \in \mathfrak{I}_{T_0}$ such that

- (i) $P(S(\mathbf{t}) \mid T_0)$ exists,
- (ii) P satisfies the principle of indifference, and
- (iii) $P(\forall x(Sx \rightarrow Gx) \mid T_0) = 1$.

Recall that in the approach of Section 7.2.4, we could not use the principle of indifference to determine the probability of success. The only difference between that approach and the approach in this section is that here, we have moved the sentence, $\forall x(Sx \rightarrow Gx)$, out of the root and into $T_\mathcal{C}$. As we discussed in Section 4.2.6, this means we are making a semantically stronger assumption about the sentence, $\forall x(Sx \rightarrow Gx)$. As we will see, this stronger assumption is enough to allow us to use the principle of indifference.

Proposition 7.2.6. *The condition \mathcal{C} is consistent and $\mathbf{P}_\mathcal{C}(S(\mathbf{t}) \mid T_0) = 1/2$.*

Proof. If P is the inductive condition constructed in the first part of the proof of Proposition 7.2.5, then $P \in \mathcal{C}$. Hence, \mathcal{C} is consistent. Let $P \in \mathcal{C}$ be arbitrary, and let $\pi = (S F)$. Then $T_0^\pi = T_0$, so that $P(S(\mathbf{t}) \mid T_0) = P(F(\mathbf{t}) \mid T_0)$, which implies $P(S(\mathbf{t}) \mid T_0) = 1/2$. Since P was arbitrary, we have $\mathbf{P}_\mathcal{C}(S(\mathbf{t}) \mid T_0) = 1/2$. \square

7.3 Examples with multiple objects

(S:basic-expls-PoI-2) In this section, we give more elementary examples involving the principle of indifference. Here, we will consider examples involving multiple objects.

7.3.1 Three balls, two colors

(S:three-balls) Imagine an urn containing three balls, each of which is black or white. The urn contains at least one white ball and at least one black ball. We will use the principle of indifference to show that every possible color combination has the same probability. The fact that there is at least one ball of each color is critical to this example. As we will see in Section 7.3.2, removing this assumption severely limits the inductive inferences that we can make.

Let $L = \{b_1, b_2, b_3, C_0, C_1\}$, where the b_k are constant symbols denoting the balls, and C_0 and C_1 are unary predicate symbols denoting the colors black and white, respectively. Let T_0 be generated by the sentences

$$\begin{aligned}\zeta_1 &: b_1 \neq b_2 \wedge b_1 \neq b_3 \wedge b_2 \neq b_3 \\ \zeta_2 &: \forall x((C_0x \vee C_1x) \wedge \neg(C_0x \wedge C_1x)) \\ \zeta_3 &: C_0b_1 \vee C_0b_2 \vee C_0b_3 \\ \zeta_4 &: C_1b_1 \vee C_1b_2 \vee C_1b_3\end{aligned}$$

As in Example 4.3.8, let $d_k(n)$ denote the k -th binary digit of n , counting digits from the right. For $n \in \{0, \dots, 7\}$, let

$$\varphi_n = C_{d_1(n)}b_1 \wedge C_{d_2(n)}b_2 \wedge C_{d_3(n)}b_3.$$

For example, 6 has the binary representation 110. Reading the digits right to left, we have 0, 1, 1. Hence, the sentence φ_6 asserts that ball b_1 is black, ball b_2 is white, and ball b_3 is white. Let \mathcal{C} be the set of $P \in \mathcal{I}_{T_0}$ such that

- (i) $P(\varphi_n | T_0)$ exists for $n \in \{0, \dots, 7\}$, and
- (ii) P satisfies the principle of indifference.

(P:three-balls) **Proposition 7.3.1.** *The condition \mathcal{C} is consistent and $\mathbf{P}_{\mathcal{C}}(\varphi_n | T_0) = 1/6$ for $n \in \{1, \dots, 6\}$.*

Proof. Let $A = \{1, 2, 3\}$. For $n \in \{0, \dots, 7\}$, define $\omega_n = (A, L^{\omega_n})$ by $b_k^{\omega_n} = k$ and $C_j^{\omega_n} = \{k \mid d_k(n) = j\}$. Note that $C_0^{\omega_n} = (C_1^{\omega_n})^c$. Also note that $\omega_n \models \varphi_m$ if and only if $m = n$. Let $\Omega = \{\omega_0, \dots, \omega_7\}$, $\Sigma = \mathfrak{P}\Omega$, and define $\mathcal{P} = (\Omega, \Sigma, \mathbb{P})$, where $\mathbb{P}\{\omega_0\} = \mathbb{P}\{\omega_7\} = 0$, and $\mathbb{P}\{\omega_n\} = 1/6$ for $1 \leq n \leq 6$. Since $\mathcal{P} \models T_0$, we may define $P = \mathbf{Th} \mathcal{P} \downarrow_{[T_0, T_0]} \in \mathcal{I}_{T_0}$. Then $P(\varphi_0 | T_0) = P(\varphi_7 | T_0) = 0$, and $P(\varphi_n | T_0) = 1/6$ for $1 \leq n \leq 6$.

We will show that P satisfies (R10). Suppose that $P(\varphi | X) = p$ and $X^\pi \in \text{ante } P$. First assume $C_0^\pi = C_0$. Let $g : A \rightarrow A$ be the bijection that satisfies $b_k^{-\pi} = b_{gk}$. Let σ be the permutation of $\{0, \dots, 7\}$ such that $d_k(\sigma n) = d_{gk}(n)$. Define $h : \mathcal{P} \rightarrow \mathcal{P}^\pi$ by $\omega_n \mapsto \omega_{\sigma n}^\pi$. Since $\sigma 0 = 0$ and $\sigma 7 = 7$, we have $\mathbb{P}^\pi\{\omega_{\sigma n}^\pi\} = \mathbb{P}\{\omega_{\sigma n}\} = \mathbb{P}\{\omega_n\}$. Hence, h induces a measure space isomorphism. Also, $\omega_{\sigma n}^\pi = \omega_{\sigma n} \circ \pi^{-1} = g \circ \omega_n$, so that $\omega_{\sigma n}^\pi \simeq \omega_n$, and h is a model isomorphism. By Theorems 7.1.7 and 5.3.24, this gives $P(\varphi^\pi | X^\pi) = p$.

Now assume $C_0^\pi = C_1$. Define $h : \mathcal{P} \rightarrow \mathcal{P}^\pi$ by $\omega_{7-n} \mapsto \omega_{\sigma n}^\pi$. As above, h induces a measure space isomorphism and $\omega_{\sigma n}^\pi = \omega_{\sigma n} \circ \pi^{-1} = g \circ \omega_{7-n}$, so that

again, $P(\varphi^\pi \mid X^\pi) = p$. This shows that P satisfies (R10), and therefore, \mathcal{C} is consistent.

Now let $P \in \mathcal{C}$ be arbitrary. Let $\pi = (b_1 \ b_2)$. Then T_0 is invariant under π and $\varphi_1^\pi = \varphi_2$, so by the principle of indifference, $P(\varphi_1 \mid T_0) = P(\varphi_2 \mid T_0)$. Similarly, using $\pi = (b_2 \ b_3)$, we have $P(\varphi_2 \mid T_0) = P(\varphi_4 \mid T_0)$. Now let $\pi = (C_0 \ C_1)$. Then $T_0^\pi = T_0$ and $\varphi_n^\pi = \varphi_{7-n}$. Thus, $P(\varphi_3 \mid T_0) = P(\varphi_4 \mid T_0)$, $P(\varphi_5 \mid T_0) = P(\varphi_2 \mid T_0)$, and $P(\varphi_6 \mid T_0) = P(\varphi_1 \mid T_0)$. It follows that for some $p \in [0, 1]$, we have $P(\varphi_n \mid T_0) = p$ for all $n \in \{1, \dots, 6\}$. But $T_0 \vdash \neg(\varphi_0 \vee \varphi_7)$, so $P(\varphi_0 \mid T_0) = P(\varphi_7 \mid T_0) = 0$. Therefore, $\sum_{n=1}^6 P(\varphi_n \mid T_0) = 1$, which implies $p = 1/6$. \square

7.3.2 Two balls, two colors

(S:two-balls) Now imagine an urn containing two balls, each of which is black or white. There are four possible color combinations. Unlike the example in Section 7.3.1, we will not be able to use the principle of indifference to show that each combination has probability $1/4$. The most we can conclude is that the probability of two whites is the same as the probability of two blacks.

Let $L = \{b_1, b_2, C_0, C_1\}$, where the b_k are constant symbols denoting the balls, and C_0 and C_1 are unary predicate symbols denoting the colors black and white, respectively. Let T_0 be generated by the sentences

$$\begin{aligned} \zeta_1 &: b_1 \neq b_2 \\ \zeta_2 &: \forall x((C_0x \vee C_1x) \wedge \neg(C_0x \wedge C_1x)) \end{aligned}$$

As in Example 4.3.8, let $d_k(n)$ denote the k -th binary digit of n , counting digits from the right. For $n \in \{0, 1, 2, 3\}$, let

$$\varphi_n = C_{d_1(n)}b_1 \wedge C_{d_2(n)}b_2.$$

For example, 2 has the binary representation 10. Reading the digits right to left, we have 0, 1. Hence, the sentence φ_2 asserts that ball b_1 is black, and ball b_2 is white. Let \mathcal{C} be the set of $P \in \mathfrak{J}_{T_0}$ such that

- (i) $P(\varphi_n \mid T_0)$ exists for $n \in \{0, 1, 2, 3\}$, and
- (ii) P satisfies the principle of indifference.

(P:two-balls) **Proposition 7.3.2.** *The condition \mathcal{C} is consistent. Moreover, for any $P \in \mathcal{C}$, we have $P(\varphi_0 \mid T_0) = P(\varphi_3 \mid T_0)$ and $P(\varphi_1 \mid T_0) = P(\varphi_2 \mid T_0)$.*

Proof. Let $A = \{1, 2\}$. For $n \in \{0, 1, 2, 3\}$, define $\omega_n = (A, L^{\omega_n})$ by $b_k^{\omega_n} = k$ and $C_j^{\omega_n} = \{k \mid d_k(n) = j\}$. Note that $C_0^{\omega_n} = (C_1^{\omega_n})^c$. Also note that $\omega_n \models \varphi_m$ if and only if $m = n$. Let $\Omega = \{\omega_0, \dots, \omega_7\}$ and $\Sigma = \mathfrak{P}\Omega$. Fix $p \in (0, 1)$ and define $\mathcal{P} = (\Omega, \Sigma, \mathbb{P})$, where $\mathbb{P}\{\omega_0\} = \mathbb{P}\{\omega_3\} = p/2$, and $\mathbb{P}\{\omega_1\} = \mathbb{P}\{\omega_2\} = (1-p)/2$. Since $\mathcal{P} \models T_0$, we may define $P = \mathbf{Th} \mathcal{P} \upharpoonright_{[T_0, \text{Th } \mathcal{P}]} \in \mathfrak{J}_{T_0}$. Then $P(\varphi_0 \mid T_0) = P(\varphi_3 \mid T_0) = p/2$, and $P(\varphi_1 \mid T_0) = P(\varphi_2 \mid T_0) = (1-p)/2$. The proof that P

satisfies (R10) follows as in the proof of Proposition 7.3.1. Hence, $P \in \mathcal{C}$ and \mathcal{C} is consistent.

Let $P \in \mathcal{C}$ be arbitrary. If $\pi = (b_1 \ b_2)$, then T_0 is invariant under π and $\varphi_1^\pi = \varphi_2$. Hence, by the principle of indifference, $P(\varphi_1 \mid T_0) = P(\varphi_2 \mid T_0)$. Similarly, if $\pi = (C_0 \ C_1)$, then $T_0^\pi = T_0$ and $\varphi_0 = \varphi_3$. Therefore, $P(\varphi_0 \mid T_0) = P(\varphi_3 \mid T_0)$. \square

Remark 7.3.3. The proof of Proposition 7.3.2 shows that \mathcal{C} is indeterminate. More specifically, for any $p \in (0, 1)$, there exists $P \in \mathcal{C}$ such that $P(\varphi_0 \mid T_0) = P(\varphi_3 \mid T_0) = p/2$, and $P(\varphi_1 \mid T_0) = P(\varphi_2 \mid T_0) = (1 - p)/2$. Therefore, $\mathbf{P}_{\mathcal{C}}(\varphi_n \mid T_0)$ does not exist for any $n \in \{0, 1, 2, 3\}$.

This example can be generalized to any finite number of black and white balls. Suppose there are N balls and let ψ_m be the sentence which asserts that exactly m of them are white. As above, the principle of indifference will be unable to tell us the probabilities of ψ_m . The most it can say is that the probability of ψ_m is the same as the probability of ψ_{N-m} .

7.3.3 Random numbers

(S:0<1) In this example, we consider a constant that could equal either 0 or 1. We will not include everything we know about the numbers 0 and 1, but we will include the fact that $0 < 1$. This creates an informational asymmetry, like the one we encountered in Section 7.2.4. As such, the principle of indifference will not provide us with the probability that this constant is equal to 0.

Let $L = \{c, \underline{0}, \underline{1}, <\}$, where c , $\underline{0}$, and $\underline{1}$ are constant symbols, and $<$ is a binary relation symbols. Let T_0 be generated by the sentences

$$\begin{aligned}\varphi_1 : \underline{0} &\neq \underline{1} \\ \varphi_2 : \underline{0} &< \underline{1} \\ \varphi_3 : c &= \underline{0} \vee c = \underline{1}\end{aligned}$$

Let $\mathcal{C} = \{P \in \mathcal{I}_{T_0} \mid P(c = \underline{0} \mid T_0) \text{ exists and } P \text{ satisfies (R10)}\}$.

(P:0<1) **Proposition 7.3.4.** *The condition \mathcal{C} is consistent and $\mathcal{C} \not\vdash (T_0, c = \underline{0}, 1/2)$.*

Proof. Let $A = \{0, 1\}$. Define $\omega_0 = (A, L^{\omega_0})$ by $\underline{0}^{\omega_0} = 0$, $\underline{1}^{\omega_0} = 1$, $<^{\omega_0} = \{(0, 1)\}$, and $c^{\omega_0} = 0$. Define ω_1 similarly, but with $c^{\omega_1} = 1$. Let $\Omega = \{\omega_0, \omega_1\}$ and $\Sigma = \mathfrak{P}\Omega$. Let $p \in (0, 1)$ and define $\mathcal{P} = (\Omega, \Sigma, \mathbb{P})$ so that $\mathbb{P}\{\omega_0\} = p$ and $\mathbb{P}\{\omega_1\} = 1 - p$. Since $\mathcal{P} \models T_0$, we may define $P = \mathbf{Th} \mathcal{P} \upharpoonright_{[T_0, Th \mathcal{P}]}$. Then $P \in \mathcal{I}_{T_0}$ and $P(c = \underline{0} \mid T_0) = p$.

We will show that P satisfies (R10). Suppose that $P(\varphi \mid X) = p$ and $X^\pi \in \text{ante } P$. Note that $<^\pi = <$ for all permutations π . First assume that $\underline{0}^\pi = \underline{1}$. Then $(\underline{0} < \underline{1})^\pi = (\underline{1} < \underline{1}^\pi)$. But $\omega \not\models (\underline{1} < s)$ for all $\omega \in \Omega$ and $s \in L$, so this contradicts Lemma 7.1.10. Hence, $\underline{0}^\pi \neq \underline{1}$. By reversing the roles of π and π^{-1} in this argument, we may also conclude that $\underline{1}^\pi \neq \underline{0}$.

It follows that if $c^\pi = c$, then π is the identity. We may therefore assume that $c^\pi = \underline{0}$ or $c^\pi = \underline{1}$. Suppose that $c^\pi = \underline{0}$, so that $\pi = (c \ \underline{0})$. We will apply

Proposition 7.1.11 with $\zeta = (\underline{0} < \underline{1})$. In this case, $B = B' = \zeta_{\Omega}^{\pi} = \{\omega_0\}$. Hence, the function $h : \mathcal{P}_B \rightarrow \mathcal{P}_{B'}$, that maps ω_i to ω_i^{π} it induces a measure space isomorphism. Moreover, $\omega_0 = \omega_0^{\pi}$, so h is a model isomorphism. Proposition 7.1.11 therefore implies $P(\varphi^{\pi} \mid X^{\pi}) = p$. A similar argument gives the same result in the case that $c^{\pi} = \underline{1}$. This shows that P satisfies (R10), so that \mathcal{C} is consistent. Since p was arbitrary, we have $\mathcal{C} \not\vdash (T_0, c = \underline{0}, 1/2)$. \square

7.3.4 Random numbers and definitions

$\langle S:0<1\text{-defn} \rangle$ In this example, we again consider a constant c that could equal 0 or 1. This time, however, we will expand our language by defining $d = 1 - c$. After making this seemingly harmless addition, we will be able to infer that c equals 0 with probability 1/2.

This result may seem counterintuitive. It feels as if introducing a defined constant should not affect the probabilities of a pre-existing constant. But this feeling is rooted in the intuition that c is the original constant, and d is defined in terms of c . However, in the expanded language, it is impossible to tell which of c and d is the original constant. They are simply two constants related by $d = 1 - c$ and $c = 1 - d$. As such, when we refer to c in the expanded language, we are equally ignorant about whether c is the original random number, or its inversion. Therefore, the probability that c equals 0 in the expanded language should be the average of its probabilities in the original language, which is 1/2. We will return to this idea at the end of the section, after formalizing the example.

Let L and T_0 be as in Section 7.3.3. That is, $L = \{c, \underline{0}, \underline{1}, <\}$ and T_0 is generated by

$$\begin{aligned}\varphi_1 : \underline{0} &\neq \underline{1} \\ \varphi_2 : \underline{0} &< \underline{1} \\ \varphi_3 : c &= \underline{0} \vee c = \underline{1}\end{aligned}$$

Let d be a constant symbol and define

$$\delta(y) : c = \underline{0} \wedge y = \underline{1} \vee c = \underline{1} \wedge y = \underline{0}$$

Then $T_0 \vdash \xi$, where $\xi = \exists! y \delta(y)$, so that $\theta = (y = d \leftrightarrow \delta(y))$ is legitimate in T_0 .

Let $L' = \{c, d, \underline{0}, \underline{1}, <\}$ and define $T'_0 = T_0 + \theta \subseteq (\mathcal{L}')^0$, so that T'_0 is a definitorial extension as in Definition 6.1.6. Let

$$\mathcal{C} = \{P' \in \mathfrak{I}_{T'_0} \mid P'(c = \underline{0} \mid T'_0) \text{ exists and } P' \text{ satisfies (R10)}\}.$$

$\langle P:0<1\text{-defn} \rangle$ **Proposition 7.3.5.** *The condition \mathcal{C} is consistent and $\mathbf{P}_{\mathcal{C}}(c = \underline{0} \mid T'_0) = 1/2$.*

Proof. Let $A = \{0, 1\}$. For $i \in \{0, 1\}$, define $\omega_i = (A, L^{\omega_i})$ by $\underline{0}^{\omega_i} = 0$, $\underline{1}^{\omega_i} = 1$, $<^{\omega_i} = \{(0, 1)\}$, and $c^{\omega_i} = i$. Let $\Omega = \{\omega_0, \omega_1\}$, $\Sigma = \mathfrak{P}\Omega$, and define $\mathcal{P} = (\Omega, \Sigma, \mathbb{P})$ so that $\mathbb{P}\{\omega_0\} = \mathbb{P}\{\omega_1\} = 1/2$. Note that \mathcal{P} is the \mathcal{L} -model defined in the proof of Proposition 7.3.4, where $p = 1/2$. Also note that $\omega \models \xi$ for all $\omega \in \Omega$. Since $\mathcal{P} \models T_0$, we may define $P = \mathbf{Th} \mathcal{P} \upharpoonright_{[T_0, Th \mathcal{P}]}$.

Let $\omega'_i = (A, (L')^{\omega'_i})$ be given by $s^{\omega'_i} = s^{\omega_i}$ for $s \in L$ and $d^{\omega'_i} = i'$, where $i' = 1 - i$. Let $\Omega' = \{\omega'_0, \omega'_1\}$, $\Gamma = \mathfrak{P}\Sigma'$, and define $\mathcal{P}' = (\Omega', \Gamma, \mathbb{Q})$ so that $\mathbb{Q}\{\omega'_0\} = \mathbb{Q}\{\omega'_1\} = 1/2$. Note that \mathcal{P}' is the \mathcal{L}' -model defined above Lemma 6.1.2. By Corollary 6.1.11, we have $P' = \mathbf{Th} \mathcal{P}' \downarrow_{[T'_0, Th \mathcal{P}']}$, where P' is the definitorial extension of P given in Theorem 6.1.10. Note that $P' \in \mathcal{J}_{T'_0}$ and $P'(c = \underline{0} \mid T'_0) = 1/2$.

We will show that P' satisfies (R10). Suppose that $P'(\varphi \mid X) = p$ and $X^\pi \in \text{ante } P'$. As in the proof of Proposition 7.3.4, we have $\underline{0}^\pi \neq \underline{1}$ and $\underline{1}^\pi \neq \underline{0}$. First assume $\pi = (c \underline{0})$. As in the proof of Proposition 7.3.4, we can use Proposition 7.1.11 with $\zeta = (\underline{0} < \underline{1})$ to conclude that $P'(\varphi^\pi \mid X^\pi) = p$. A similar argument covers the cases where π is $(c \underline{1})$, $(d \underline{0})$, or $(d \underline{1})$. If $\pi = (c d)$, then $(\omega'_i)^\pi = \omega'_i$, so that $(\mathcal{P}')^\pi \simeq \mathcal{P}'$, which also gives $P'(\varphi^\pi \mid X^\pi) = p$. This covers the case that π is a transposition.

Now suppose π is a 3-cycle. Assume $\pi = (\underline{0} c d)$. Then $B = \zeta_\Omega^\pi = \{\omega'_0\}$ and $B' = \zeta_{\Omega'}^\pi = \{\omega'_1\}$. Since $(\omega'_1)^\pi = \omega'_0$, it follows from Proposition 7.1.11 that $P'(\varphi^\pi \mid X^\pi) = p$. A similar argument covers the cases where π is $(\underline{0} d c)$, $(\underline{1} c d)$, or $(\underline{1} d c)$.

Finally, suppose π affects every constant symbol. If $\pi = (c \underline{0})(d \underline{1})$, then $(\omega'_i)^\pi = \omega'_i$, so that $(\mathcal{P}')^\pi = \mathcal{P}'$, which gives $P'(\varphi^\pi \mid X^\pi) = p$. A similar argument covers $\pi = (c \underline{1})(d \underline{0})$. The remaining possibility is that π is a 4-cycle. Assume $\pi = (\underline{0} c \underline{1} d)$. As above, we have $B = \{\omega'_0\}$ and $B' = \{\omega'_1\}$, so that Proposition 7.1.11 gives $P'(\varphi^\pi \mid X^\pi) = p$. A similar argument covers the case $\pi = (\underline{0} d \underline{1} c)$. Altogether, this shows that P' satisfies (R10), so that \mathcal{C} is consistent.

Now let $P' \in \mathcal{C}$ be arbitrary. Let $\pi = (c d)$. Then T'_0 is invariant under π . Hence, $P'(c = \underline{0} \mid T'_0) = P'(d = \underline{0} \mid T'_0)$. But $d = \underline{0} \equiv_{T'_0} c = \underline{1}$. Therefore, $P'(c = \underline{0} \mid T'_0) = P'(c = \underline{1} \mid T'_0)$, which gives $P'(c = \underline{0} \mid T'_0) = 1/2$. \square

This example can be generalized to a constant c that could equal 0, 1, or 2. That is, let $L = \{c, \underline{0}, \underline{1}, \underline{2}\}$ and let T_0 be generated by

$$\begin{aligned} \varphi_1 : \underline{0} \neq \underline{1} \wedge \underline{0} \neq \underline{2} \wedge \underline{1} \neq \underline{2} \\ \varphi_2 : \underline{0} < \underline{1} \wedge \underline{0} < \underline{2} \wedge \underline{1} < \underline{2} \\ \varphi_3 : c = \underline{0} \vee c = \underline{1} \vee c = \underline{2} \end{aligned}$$

As in Section 7.3.3, we cannot use the principle of indifference to determine $P(c = \underline{n} \mid T_0)$. But we can create a definitorial expansion T'_0 using

$$\theta_d = (y = d \leftrightarrow c = \underline{0} \wedge y = \underline{1} \vee c = \underline{1} \wedge y = \underline{0} \vee c = \underline{2} \wedge y = \underline{2}).$$

Informally, $d = f(c)$, where f interchanges 0 and 1. As above, we could then use the principle of indifference to conclude that $P'(c = \underline{0} \mid T'_0) = P'(c = \underline{1} \mid T'_0)$.

It is tempting to think we can iterate this process. That is, suppose we create a definitorial expansion T''_0 of T'_0 using

$$\theta_e = (y = e \leftrightarrow c = \underline{0} \wedge y = \underline{2} \vee c = \underline{1} \wedge y = \underline{1} \vee c = \underline{2} \wedge y = \underline{0}).$$

Informally, $e = g(c)$, where g interchanges 0 and 2. We might now expect that the principle of indifference gives $P''(c = \underline{2} \mid T_0'') = 1/3$. But it does not. In fact, even our previous inference is no longer valid. That is, we can no longer even conclude that $P''(c = \underline{0} \mid T_0'') = P''(c = \underline{1} \mid T_0'')$. This is because T_0'' is no longer invariant under $\pi = (c\ d)$. In particular, $T_0'' \vdash d = \underline{0} \rightarrow e = \underline{1}$, but $T_0'' \not\vdash c = \underline{0} \rightarrow e = \underline{1}$.

Remark 7.3.6. This example shows that definitorial extensions do not preserve the principle of indifference. It is possible for P to satisfy (R10), but for its definitorial extension P' to not satisfy it. The converse is also possible. In other words, Theorem 6.1.10 is another theorem that would fail if we included (R10) in the definition of an inductive theory.

The juxtaposition of Propositions 7.3.4 and 7.3.5 may seem counterintuitive. On the one hand, in \mathcal{L} , we cannot infer the probability that $c = \underline{0}$. On the other hand, by passing to \mathcal{L}' , whose only difference is that it includes the defined constant d , we are suddenly able to infer that $c = \underline{0}$ with probability $1/2$. It seems that we must have added some new information by passing to \mathcal{L}' . But clearly we did not. It is the nature of a definitorial extension that it adds no new logical information. The explanation is not that we have added new information. Rather, we have altered the very meaning of c , in the way described in Section 6.1.6.

It is tempting to think that a constant symbol such as c stands for some object. From that point of view, it must stand for the same object in both \mathcal{L} and \mathcal{L}' . And in that case, it makes no sense to say that we altered the meaning of c . But syntactically, c does not stand for anything. Standing for an object is a semantic notion. What c stands for is relative to the model we are using, and even then, it can vary from structure to structure within that model. Syntactically speaking, c is not denoting an object. Rather, it is a primitive symbol that gains its meaning from the deductive and inductive facts that use it.

In \mathcal{L}' , we have changed those facts from T_0 to T_0' . The meaning of c in T_0' is not necessarily the same as in T_0 . To see this more clearly, simply rename c and d in \mathcal{L}' to c' and d' . It is then no longer surprising that $P'(c' = \underline{0} \mid T_0') = 1/2$. After all, in T_0' , it is impossible to tell which of c' and d' is the original constant from \mathcal{L} , and which of them is defined in terms of that original constant. We are indifferent about those two possibilities. Therefore, $P'(c' = \underline{0} \mid T_0')$ should be the average of $P(c = \underline{0} \mid T_0)$ and $P(c = \underline{1} \mid T_0)$, which is $1/2$. By analogy with measure-theoretic probability, it is as if we started with a $\{0, 1\}$ -valued random variable X , defined $Y = 1 - X$, and then let (X', Y') be a random permutation of (X, Y) . In that case, if $(X', Y') = (X, Y)$ and $(X', Y') = (Y, X)$ are equally likely, then $P(X' = 0) = 1/2$, regardless of the distribution of X .

7.4 Indifference and exchangeability

(S:indiff-exch) For the remainder of this chapter, we will look at how the principle of indifference relates to real inductive theories. For simplicity, and to match the intuition of measure-theoretic probability models, we will take the approach presented in Theorem 6.4.6. Namely, we will operate under the standing assumption that ZFC is strictly satisfiable. This assumption will be in effect for the remainder of this chapter.

In this short section, we show that, in the context of a measure-theoretic probability model, exchangeability is a special case of the principle of indifference. This result is given below in Theorem 7.4.2. In order to prove it, we first establish Lemma 7.4.1, which we be useful several times throughout the remainder of this chapter.

7.4.1 Permutations of real inductive theories

Let $P \subseteq \mathcal{L}^{\text{IS}}$ be a real inductive theory in ZFC. Let π be an L -permutation such that $\in^\pi = \in$. Define $\pi' : L \rightarrow L$ by

$$s^{\pi'} = \begin{cases} s & \text{if } s \in L_{\text{ZFC}}, \\ s^\pi & \text{if } s \notin L_{\text{ZFC}} \text{ and } s^\pi \notin L_{\text{ZFC}}, \text{ and} \\ (s^\pi)^\pi & \text{if } s \notin L_{\text{ZFC}} \text{ and } s^\pi \in L_{\text{ZFC}}. \end{cases}$$

(L:fixed-perm) **Lemma 7.4.1.** *Suppose $X, X^\pi \in \text{ante } P$. Then the function π' is an L -permutation that fixes L_{ZFC} . Moreover, if $X^{\pi'} \in \text{ante } P$ and π' satisfies (R10), then π satisfies (R10).*

Proof. Assume $X, X^\pi \in \text{ante } P$. By construction, the function π' fixes L_{ZFC} and preserves the type and arity of extralogical symbols. The fact that π' is a bijection is a consequence of the following:

$$\text{if } s \in L_{\text{ZFC}} \text{ and } s^\pi \neq s, \text{ then } s^{-\pi} \notin L_{\text{ZFC}} \text{ and } s^\pi \notin L_{\text{ZFC}}. \quad (7.4.1) \text{fixed-perm}$$

To see this, let $s \in L_{\text{ZFC}}$ with $s^\pi \neq s$. Since $\in^\pi = \in$, we have $s \neq \in$. Hence, s is an explicitly defined constant symbol. Let $\delta(y)$ be its defining formula, and let $\delta^{\text{rd}}(y)$ be its reduction to $\mathcal{L}\{\in\}$, so that the only extralogical symbol in $\delta^{\text{rd}}(y)$ is \in . Let $\zeta = \forall y(y = s \leftrightarrow \delta^{\text{rd}}(y))$, so that $\text{ZFC} \vdash \zeta$. By Lemma 7.1.10, if $\mathcal{P} \models P$, then $\psi_\Omega \subseteq \zeta_\Omega^{-\pi}$ a.s., which implies $\mathcal{P} \models (X, \zeta^{-\pi}, 1)$. Hence, $P(\zeta^{-\pi} \mid X) = 1$. It follows that $P(\zeta \wedge \zeta^{-\pi} \mid X) = 1$. But $\zeta^{-\pi} = \forall y(y = s^{-\pi} \leftrightarrow \delta^{\text{rd}}(y))$, so that $\text{ZFC} \vdash \zeta \wedge \zeta^{-\pi} \leftrightarrow s = s^{-\pi}$. Therefore $P(s = s^{-\pi} \mid X) = 1$. A similar argument shows that $P(\zeta^\pi \mid X^\pi) = 1$ and $P(s = s^\pi \mid X^\pi) = 1$. Finally, note that if $s, s' \in L_{\text{ZFC}} \setminus \{\in\}$, then $\text{ZFC} \vdash s \neq s'$. Hence, $s^{-\pi} \notin L_{\text{ZFC}}$ and $s^\pi \notin L_{\text{ZFC}}$.

Now assume $X^{\pi'} \in \text{ante } P$ and π' satisfies (R10). Using the facts that $\in^\pi = \in$, every $s \in L_{\text{ZFC}} \setminus \{\in\}$ is a constant symbol, and $P(s = s^\pi \mid X^\pi) = 1$ for all $s \in L_{\text{ZFC}} \setminus \{\in\}$, it follows by term induction and formula induction that

$$P(\theta^{\pi'} \leftrightarrow \theta^\pi \mid X^\pi) = 1 \text{ for all } \theta \in \mathcal{L}^0. \quad (7.4.2) \text{fixed-perm-1}$$

A similar argument shows that $P(\theta^{-\pi'} \leftrightarrow \theta^{-\pi} \mid X) = 1$ for all $\theta \in \mathcal{L}^0$. Applying this to θ^π gives $P((\theta^\pi)^{-\pi'} \leftrightarrow \theta \mid X) = 1$. But π' satisfies (R10), so this gives

$$P(\theta^\pi \leftrightarrow \theta^{\pi'} \mid X^{\pi'}) = 1 \text{ for all } \theta \in \mathcal{L}^0. \quad (7.4.3) \text{ fixed-perm-2}$$

To show that π satisfies (R10), suppose that $P(\varphi \mid X) = p$. Since π' satisfies (R10), we have $P(\varphi^{\pi'} \mid X^{\pi'}) = p$. Also, (7.4.3) gives $P(\varphi^\pi \leftrightarrow \varphi^{\pi'} \mid X^{\pi'}) = 1$. Hence, Proposition 3.2.14 implies $P(\varphi^\pi \mid X^{\pi'}) = p$.

Let $\theta \in X^\pi$. Then $(\theta^{-\pi})^{\pi'} \in X^{\pi'}$, so that $P((\theta^{-\pi})^{\pi'} \mid X^{\pi'}) = 1$. As it was for φ above, this gives $P((\theta^{-\pi})^\pi \mid X^{\pi'}) = 1$. But $(\theta^{-\pi})^\pi = \theta$. Hence, $P(\theta \mid X^{\pi'}) = 1$ for all $\theta \in X^\pi$. By the rule of deductive extension, $P(\cdot \mid X^{\pi'}, X^\pi) = P(\cdot \mid X^{\pi'})$. In particular, $P(\varphi^\pi \mid X^{\pi'}, X^\pi) = p$. On the other hand, the analogous argument using (7.4.2) shows that $P(\theta \mid X^\pi) = 1$ for all $\theta \in X^{\pi'}$. The rule of deductive extension therefore also shows that $P(\cdot \mid X^{\pi'}, X^\pi) = P(\cdot \mid X^\pi)$. Hence, $P(\varphi^\pi \mid X^\pi) = p$. \square

7.4.2 Exchangeability

Let $\langle X_i \mid i \in I \rangle$ be a collection of real-valued random variables defined on a probability space, (S, Γ, ν) . We say that $\langle X_i \mid i \in I \rangle$ are *exchangeable* if the distribution of $(X_{i(1)}, \dots, X_{i(n)})$ is unchanged by a finite permutation of I . More specifically, let $\sigma : I \rightarrow I$ be a bijection with $\sigma(i) = i$ for all but finitely many i . Then

$$\nu \bigcap_{k=1}^n \{X_{i(k)} \in V_k\} = \nu \bigcap_{k=1}^n \{X_{\sigma(i(k))} \in V_k\}, \quad (7.4.4) \text{ exchangeable}$$

for all choices of $n \in \mathbb{N}$, $i(k) \in I$, and $V_k \in \mathcal{B}(\mathbb{R})$.

Let $\langle X_i \mid i \in I \rangle$ be real-valued random variables on (S, Γ, ν) . Without loss of generality, we may assume $S = \mathbb{R}^I$, $\Gamma = \bigotimes_{i \in I} \mathcal{B}(\mathbb{R})$, and $X_i(x) = x_i$. Recall our standing assumption that ZFC is strictly satisfiable. Let $\mathcal{P} = (\Omega, \Sigma, \mathbb{P})$ be the model constructed in the proof of Theorem 6.4.6, and let $P = \mathbf{Th} \mathcal{P} \upharpoonright_{[\text{ZFC}, \text{Th } \mathcal{P}]}$.

(T:exchangeable) Theorem 7.4.2. *With notation as above, the inductive theory P satisfies the principle of indifference if and only if $\langle X_i \mid i \in I \rangle$ are exchangeable.*

Proof. Assume P satisfies the principle of indifference. Let $\sigma : I \rightarrow I$ be a bijection with $\sigma(i) = i$ for all but finitely many i . Let π be the signature permutation that fixes everything in L_{ZFC} , and maps \underline{X}_i to $\underline{X}_{\sigma(i)}$. Then $\psi^\pi = \psi$ for all $\psi \in \mathcal{L}_{\text{ZFC}}$. In particular, $\text{ZFC}^\pi = \text{ZFC}$. Let $\varphi = \bigwedge_{k=1}^n \underline{X}_{i(k)} \in \underline{V}_k$. Then $\varphi^\pi = \bigwedge_{k=1}^n \underline{X}_{\sigma(i(k))} \in \underline{V}_k$. By the principle of indifference,

$$P(\bigwedge_{k=1}^n \underline{X}_{i(k)} \in \underline{V}_k \mid \text{ZFC}) = P(\bigwedge_{k=1}^n \underline{X}_{\sigma(i(k))} \in \underline{V}_k \mid \text{ZFC}).$$

Therefore, by (6.4.2), we have (7.4.4).

Now assume X is exchangeable. Let π be an L -permutation and suppose $P(\varphi \mid Y) = p$. Since \in is the only binary operation symbol in L , we have $\in^\pi = \in$. By Lemma 7.4.1, it suffices to assume π fixes L_{ZFC} , and to show that both $Y^\pi \in \text{ante } P$ and $P(\varphi^\pi \mid Y^\pi) = p$.

Since π fixes L_{ZFC} , it only affects $C = \{\underline{X}_i \mid i \in I\}$. Let $\sigma : I \rightarrow I$ be the bijection defined by $\underline{X}_i^\pi = \underline{X}_{\sigma(i)}$ and define the bijection $g : S \rightarrow S$ by $(gx)_{\sigma(i)} = x_i$. Note that both g and g^{-1} are measurable. We claim that $\nu = \nu \circ g^{-1}$. To verify this, it suffices to check that $\nu B = \nu g^{-1}B$, when B is a cylinder set of the form

$$B = \bigcap_{k=1}^n \{x \in S \mid x_{\sigma(i(k))} \in V_k\}.$$

In this case, we have

$$\begin{aligned} g^{-1}B &= \bigcap_{k=1}^n \{x \in S \mid (gx)_{\sigma(i(k))} \in V_k\} \\ &= \bigcap_{k=1}^n \{x \in S \mid x_{i(k)} \in V_k\}, \end{aligned}$$

so that $\nu B = \nu g^{-1}B$ follows from (7.4.4). Hence, $\nu = \nu \circ g^{-1}$, so that g is a pointwise isomorphism from (S, Γ, ν) to itself.

Let $h : S \rightarrow \Omega$ be the function in the proof of Theorem 6.4.6 that maps $x \in S$ to $\omega^x \in \Omega$. Let $\mathcal{P}^\pi = (\Omega^\pi, \Sigma^\pi, \mathbb{Q})$ and let $h_\pi : \Omega \rightarrow \Omega^\pi$ be the function that maps ω to ω^π . Then \mathcal{P}^π is the measure space image of (S, Γ, ν) under $h_\pi \circ h$. But $h_\pi \circ h = h \circ g$ and g is a pointwise isomorphism, so \mathcal{P}^π is the measure space image of (S, Γ, ν) under h . By the definition of \mathcal{P} , this means $\mathcal{P}^\pi = \mathcal{P}$.

Now, since $P(\varphi \mid Y) = p$ exists, we have $\mathcal{P} \models (Y, \varphi, p)$. By Theorem 7.1.7 and $\mathcal{P} = \mathcal{P}^\pi$, it follows that $\mathcal{P} \models (Y^\pi, \varphi^\pi, p)$. Therefore, $Y^\pi \in \text{ante } P$ $P(\varphi^\pi \mid Y^\pi) = p$. \square

7.5 Examples on an interval

$\langle \text{S:indiff-interval} \rangle$ In this section, we present examples of the principle of indifference that involve an interval on the real line.

7.5.1 The interval $[0, 1]$

$\langle \text{S:[0,1]} \rangle$ In our first example, we have a real number c , about which we know only that $c \in [0, 1]$. We then ask what the principle of indifference has to say about the distribution of c . At first glance, we might expect the principle to assign c a uniform distribution, based on the fact that we are somehow “equally ignorant” about where c lies in the interval $[0, 1]$. A little further thought, however, quickly reveals that this cannot be the case. The principle of indifference requires an informational symmetry, encoded in the permutation π . In the coin flip of Section 7.2.2, for example, we obtained the probability $1/2$ by interchanging the symbols for heads and tails. We could do this because our background information, T_0 , was symmetric with respect to this interchange.

In this case, however, our background information, T_0 , will contain ZFC, and the individual numbers in $[0, 1]$ are most certainly not interchangeable with respect to ZFC. For instance, in ZFC, we know that $\underline{0}$ is the additive identity and $\underline{1}$ is the multiplicative identity. That sentence is no longer true if we interchange

$\underline{0}$ and $\underline{1}$. To say that $c = \underline{0}$ is a qualitatively different assertion than to say that $c = \underline{1}$.

In this way, the present example is less like the coin flip of Section 7.2.2 and more like the random number (either $\underline{0}$ or $\underline{1}$) in Section 7.3.3. In that example, the principle of indifference did not narrow down the probabilities at all. Every possible value for the probability of $\underline{0}$ was consistent with the principle of indifference. Similarly, here, every possible distribution on c will be consistent with it.

To state this result, let L_{ZFC} be the logical signature of ZFC, given by (6.4.1). Let c be a constant symbol not in L_{ZFC} , let $L = L_{\text{ZFC}} \cup \{c\}$, and define $T_0 = \text{ZFC} + c \in [0, 1]$. Let \mathcal{C} be the inductive condition consisting of all inductive theories $P \subseteq \mathcal{L}^{\text{IS}}$ with root T_0 such that

- (i) c is Borel given T_0 , and
- (ii) P satisfies the principle of indifference.

Let ν be an arbitrary Borel probability measure on \mathbb{R} . Fix an L_{ZFC} -structure ω_0 such that $\omega_0 \models \text{ZFC}$. For each $r \in \mathbb{R}$, define $\omega = \omega^r$ to be the L -expansion of ω_0 given by $c^\omega = r^{\omega_0}$. Let $\Omega = \{\omega^r \mid r \in \mathbb{R}\}$ and let $h : \mathbb{R} \rightarrow \Omega$ denote the map $r \mapsto \omega^r$. Let $\mathcal{P}_\nu = (\Omega, \Sigma, \mathbb{P})$ be the measure space image of $(\mathbb{R}, \mathcal{B}(\mathbb{R}), \nu)$ under the function h .

If $\nu[0, 1] = 1$, then $\mathcal{P}_\nu \models T_0$, so may define the complete inductive theory $P_\nu = \mathbf{Th} \mathcal{P}_\nu \upharpoonright_{[T_0, \text{Th} \mathcal{P}_\nu]}$. By Theorem 6.4.6, we have $P(c \in \underline{V} \mid T_0) = \nu V$ for all $V \in \mathcal{B}(\mathbb{R})$. In particular, c is Borel given T_0 under P_ν .

Proposition 7.5.1. *For any such ν , we have $P_\nu \in \mathcal{C}$.*

Proof. Suppose $P_\nu(\varphi \mid X)$ exists and let π be an L -permutation. As in the proof of Theorem 7.4.2, we may assume the permutation π fixes all of L_{ZFC} . Since there is only one symbol in $L \setminus L_{\text{ZFC}}$, the permutation π is the identity. Therefore, it is trivially the case that $P_\nu(\varphi^\pi \mid X^\pi) = P_\nu(\varphi \mid X)$. \square

7.5.2 A point on a rod

(S:pt-on-rod) Introduction

In this example, we consider a rigid rod of some unspecified length, and we let c be a point that lies somewhere along the length of the rod. We then ask what, if anything, the principle of indifference has to say about the distribution of the point c .

A common approach to a problem like this would be to replace the rod with the interval $[0, 1]$. If we do that, then we have our answer, according to the preceding example: the principle of indifference says nothing. But we should ask ourselves why we feel justified in replacing the rod with $[0, 1]$. When we do so, we are introducing qualitative differences between points on the rod that were not there originally. In other words, we are making assumptions that are not indicated by the statement of the problem.

Instead of *replacing* the rod with $[0, 1]$, we should *represent* it by $[0, 1]$. For instance, we could replace the rod with a smooth manifold with boundary, and give it a Riemannian metric that indicates it has no curvature. If we did that, then we would see informational symmetries that are not there when we simply replace the rod with $[0, 1]$.

Taking this approach would require us to formulate, in ZFC, a number of new and complicated definitions. While this is certainly doable, we will avoid these complications by simply replacing our rod with a subset M of the real line, and assuming that M can be parameterized by some affine linear function on $[0, 1]$. We then let c be an element of M .

To talk about the distribution of c , we need to have names for subsets of M . For this, we will let c_0 and c_1 be the endpoints of M , arbitrarily labeled, and name the Borel subsets of M relative to c_0 . That is, if B is a Borel subset of $[0, 1]$, then B_* will be the image of B when $[0, 1]$ is mapped to M in a way that sends 0 to c_0 and 1 to c_1 .

Notation in ZFC

To make this precise, we first establish some new shorthand in \mathcal{L}_{ZFC} , the language of ZFC. Let $\delta(u, v, y) \in \mathcal{L}_{\text{ZFC}}$ be given by

$$\begin{aligned} \delta(u, v, y) &= (u \notin \mathbb{R} \vee v \notin \mathbb{R}) \wedge y = \emptyset \\ &\vee u \in \mathbb{R} \wedge v \in \mathbb{R} \wedge y \in \mathbb{R}^{\mathbb{R}} \wedge (\forall x \in \mathbb{R})(y(x) = u \cdot x + v). \end{aligned}$$

Then $\text{ZFC} \vdash \forall uv \exists! y \delta(u, v, y)$. Hence, we could explicitly define the function symbol F by $y = Fuv \leftrightarrow \delta(u, v, y)$, and then let f_{uv} be shorthand for the term Fuv . We do not, in fact, add the symbol F to our extralogical signature, but instead regard both F and f_{uv} as shorthand. We also adopt the shorthand, $\text{dom}(f)$ and $f^{\text{Img}}(z)$, given by

$$\begin{aligned} \text{dom}(f) &= \{x \in \mathbb{R} \mid (\exists y \in \mathbb{R})(x, y) \in f\}, \\ f^{\text{Img}}(z) &= \{f(x) \mid x \in z \cap \text{dom}(f)\}. \end{aligned}$$

Extralogical symbols and assumptions

To talk about our rod, we add new extralogical symbols to L_{ZFC} , the signature of ZFC. Let $\mathcal{E} = \{B \in \mathcal{B}([0, 1]) \mid \emptyset \subset B \subset [0, 1]\}$ and let

$$C = \{M, c, c_0, c_1\} \cup \{B_* \mid B \in \mathcal{E}\}$$

be a set of distinct constant symbols not in L_{ZFC} . Let $L = L_{\text{ZFC}}C$.

In the language \mathcal{L} , we want to build T_0 , the assumptions we will be making in our setup of the problem. Using $\varphi_M(u, v) = (u \neq \underline{0} \wedge M = f_{uv}^{\text{Img}}(\underline{[0, 1]}))$, define the following sentences in \mathcal{L} :

$$\begin{aligned} \varphi_1 &: (\exists uv \in \mathbb{R}) \varphi_M(u, v) \\ \varphi_2 &: c \in M \\ \varphi_3 &: (\exists uv \in \mathbb{R})(\varphi_M(u, v) \wedge c_0 = v \wedge c_1 = u + v) \end{aligned}$$

Also, for any $B \in \mathcal{E}$, define

$$\varphi_B(x) : (\exists uv \in \mathbb{R})(\varphi_M(u, v) \wedge c_0 = v \wedge x = f_{uv}^{\text{Img}}(B))$$

Then let $T_0 = \text{ZFC} + \{\varphi_1, \varphi_2, \varphi_3\} \cup \{\varphi_B(B_*) \mid B \in \mathcal{E}\}$.

The sentence φ_1 says that M is a “rod.” That is, M can be parameterized by the interval $[0, 1]$. Another way to think of this is that M is a closed interval, but its location, length, and orientation are left unspecified.

The sentence φ_2 says that c is an element of M . This is the only information about c contained in T_0 .

The sentence φ_3 says that c_0 and c_1 are the two distinct endpoints of M . But it does not specify which is the left endpoint and which is the right. Thinking of M as a rod, it does not even make sense to ask which end is left and which is right, for the orientation of the rod is arbitrary. It therefore makes sense that we would not include such information in T_0 .

The sentence $\varphi_B(B_*)$ defines the symbol B_* so that B_* is the representative of B in M , relative to c_0 . For instance, if $B = [0, 1/2]$, then B_* is the subset of M that extends from c_0 to the midpoint of M .

Inductive hypotheses and conclusion

Now let \mathcal{C} be the inductive condition consisting of all inductive theories $P \subseteq \mathcal{L}^{\text{IS}}$ with root T_0 such that

- (i) $P(c \in B_* \mid T_0)$ exists for all $B \in \mathcal{E}$, and
- (ii) P satisfies the principle of indifference.

$\langle \text{T:stick} \rangle$ **Theorem 7.5.2.** *The condition \mathcal{C} is consistent. Moreover, if $P \in \mathcal{C}$, then*

$$P(c \in B_* \mid T_0) = P(c \in (\rho B)_* \mid T_0), \quad (7.5.1) \text{stick}$$

for all $B \in \mathcal{B}([0, 1])$, where $\rho : [0, 1] \rightarrow [0, 1]$ is given by $\rho r = 1 - r$.

The proof of Theorem 7.5.2 will be given at the end of this subsection. Proving consistency is what will require the most work. This consistency proof will show, in fact, that (7.5.1) is the only restriction on the distribution of c . In particular, the principle of indifference does not require the distribution to be uniform. This is because the only points on the rod that are interchangeable with respect to T_0 are points that are equidistant from the ends. To say that c lies twice as far from one endpoint than the other is a qualitatively different statement than saying that c lies at the midpoint. These two locations are not symmetric with respect to everything we know about the rod.

A model for the rod

Our proof of consistency will utilize an \mathcal{L} -model, $\mathcal{P} = (\Omega, \Sigma, \mathbb{P})$, which we construct as follows. Let ν_0 be a probability measure on $(\mathbb{R}, \mathcal{B}(\mathbb{R}))$ such that $\nu_0[0, 1] = 1$ and ν_0 is continuous. That is, $\nu_0\{r\} = 0$ for all $r \in \mathbb{R}$. Let

$S = \mathbb{R} \times \mathbb{R} \times \{0, 1\}$ and $\Gamma = \mathcal{B}(\mathbb{R}) \otimes \mathcal{B}(\mathbb{R}) \otimes \mathfrak{P}\{0, 1\}$, and define the probability measure ν on (S, Γ) by

$$\nu B \times B' \times \{n\} = (1/2)(\nu B)(\nu B').$$

Let ω_0 be an L_{ZFC} -structure such that $\omega_0 \equiv \text{ZFC}$. If $x = (r, t, n) \in S$, then let $\omega = \omega^x$ be the L -expansion of ω_0 defined by $M^\omega = [t, t+1]^{\omega_0}$, $c^\omega = \underline{t} + r^{\omega_0}$,

$$c_0^\omega = \begin{cases} \underline{t}^{\omega_0} & \text{if } n = 0, \\ \underline{t+1}^{\omega_0} & \text{if } n = 1, \end{cases} \quad c_1^\omega = \begin{cases} \underline{t+1}^{\omega_0} & \text{if } n = 0, \\ \underline{t}^{\omega_0} & \text{if } n = 1, \end{cases}$$

and

$$B_*^\omega = \begin{cases} \tau_t B^{\omega_0} & \text{if } n = 0, \\ \tau_t \rho B^{\omega_0} & \text{if } n = 1, \end{cases}$$

where $\tau_t : \mathbb{R} \rightarrow \mathbb{R}$ and $\rho : \mathbb{R} \rightarrow \mathbb{R}$ are defined by $\tau_t r = t + r$ and $\rho r = 1 - r$.

Let $\Omega = \{\omega^x \mid x \in S\}$, let $h : S \rightarrow \Omega$ denote the function $x \mapsto \omega^x$, and let \mathcal{P} be the measure space image of (S, Γ, ν) under h .

(L:stick-1) **Lemma 7.5.3.** *With notation as above, we have $\mathcal{P} \models T_0$.*

Proof. Since each $\omega \in \Omega$ is an extension of ω_0 , we have $\mathcal{P} \models \text{ZFC}$. Note that $\omega \equiv \varphi_M[\underline{1}^{\omega_0}, \underline{t}^{\omega_0}]$ for all $\omega \in \Omega$. Hence, $(\varphi_1)_\Omega = \Omega$, so that $\mathcal{P} \models \varphi_1$. Also, $h^{-1}(\varphi_2)_\Omega = [0, 1] \times \mathbb{R} \times \{0, 1\}$, so that $\overline{\mathbb{P}}(\varphi_2)_\Omega = \nu_0[0, 1] = 1$. Thus, $\mathcal{P} \models \varphi_2$.

Note that

$$\omega_0 \equiv (x = f_{uv}^{\text{Img}}(B))[\tau_t B^{\omega_0}, \underline{1}^{\omega_0}, \underline{t}^{\omega_0}], \quad (7.5.2) \text{ \texttt{stick-1}}$$

$$\omega_0 \equiv (x = f_{uv}^{\text{Img}}(B))[\tau_t \rho B^{\omega_0}, -\underline{1}^{\omega_0}, \underline{t+1}^{\omega_0}], \quad (7.5.3) \text{ \texttt{stick-2}}$$

for all $B \in \mathcal{E}$. Fix $B \in \mathcal{E}$. Let $x = (r, t, 0)$ and $\omega = \omega^x$. Then $\omega \equiv \varphi_M[\underline{1}^{\omega_0}, \underline{t}^{\omega_0}]$, $\omega \equiv c_0 = \underline{t}$, and $\omega \equiv c_1 = \underline{1} + \underline{t}$. Also, by (7.5.2) and the definition of B_*^ω , we have $\omega \equiv (B_* = f_{uv}^{\text{Img}}(B))[\underline{1}^{\omega_0}, \underline{t}^{\omega_0}]$. It therefore follows that $\omega \equiv \varphi_3$ and $\omega \equiv \varphi_B(B_*)$. Similarly, if $x = (r, t, 1)$ and $\omega = \omega^x$, then $\omega \equiv \varphi_M[-\underline{1}^{\omega_0}, \underline{t+1}^{\omega_0}]$, $\omega \equiv c_0 = \underline{t+1}$, and $\omega \equiv c_1 = -\underline{1} + \underline{t+1}$. This time using (7.5.3) and the definition of B_*^ω , we have $\omega \equiv (B_* = f_{uv}^{\text{Img}}(B))[-\underline{1}^{\omega_0}, \underline{t+1}^{\omega_0}]$. Therefore, in this case also, we obtain $\omega \equiv \varphi_3$ and $\omega \equiv \varphi_B(B_*)$. \square

Narrowing down the permutations

By Lemma 7.5.3, we may define the inductive theory $P = \mathbf{Th} \mathcal{P} \upharpoonright_{[T_0, Th \mathcal{P}]}$. We will show that \mathcal{C} is consistent by showing that $P \in \mathcal{C}$. The difficult part of proving this is showing that P satisfies the principle of indifference. In preparation for this, we first prove two lemmas that narrow down the possible permutations that need to be checked.

(L:stick-2) **Lemma 7.5.4.** *If $P(\varphi \mid X) = p$ and $X^\pi \in \text{ante } P$, then $s^\pi = s$ for all $s \in L_{\text{ZFC}}$ and $\pi\{c, c_0, c_1\} = \{c, c_0, c_1\}$.*

Proof. Assume $M^\pi \in L_{\text{ZFC}}$. Let $\mathbf{s} = M^\pi$. By the argument following (7.4.1), we have $P(\mathbf{s} = M \mid X) = 1$. Note that

$$h^{-1}(\mathbf{s} = M)_\Omega = \{(r, t, n) \in S \mid \mathbf{s}^{\omega_0} = \underline{[t, t+1]}^{\omega_0}\}.$$

Since $t \neq t'$ implies $\text{ZFC} \vdash [t, t+1] \neq [t', t'+1]$, and since $\omega_0 \models \text{ZFC}$, there can be at most one $t \in \mathbb{R}$ such that $\mathbf{s}^{\omega_0} = \underline{[t, t+1]}^{\omega_0}$. Hence, we may choose $t_0 \in \mathbb{R}$ such that $h^{-1}(\mathbf{s} = M)_\Omega \subseteq \mathbb{R} \times \{t_0\} \times \{0, 1\}$. Since $\nu_0\{t_0\} = 0$, this implies $\overline{\mathbb{P}}(\mathbf{s} = M)_\Omega = \overline{\nu} h^{-1}(\mathbf{s} = M)_\Omega = 0$. Therefore, $P(\mathbf{s} = M \mid T_0) = 0$. By (3.2.5) and deductive transitivity, $P(\mathbf{s} = M \mid X) = 0$, a contradiction. This shows that $M^\pi \notin L_{\text{ZFC}}$. Similar arguments show that $\mathbf{s}^\pi \notin L_{\text{ZFC}}$ for all $\mathbf{s} \in C$. That is, if $\mathbf{s} \notin L_{\text{ZFC}}$, then $\mathbf{s}^\pi \notin L_{\text{ZFC}}$. By contraposition, if $\mathbf{s} \in L_{\text{ZFC}}$, then $\mathbf{s}^{-\pi} \in L_{\text{ZFC}}$. Hence, by (7.4.1), we have $\mathbf{s}^\pi = \mathbf{s}$ for all $\mathbf{s} \in L_{\text{ZFC}}$.

Now assume $c^\pi \notin \{c, c_0, c_1\}$. By the above, $c^\pi \notin L_{\text{ZFC}}$. Hence, either $c^\pi = M$ or $c^\pi = B_*$ for some $B \in \mathcal{E}$. In either case, $\omega \models c^\pi \notin \mathbb{R}$ for all $\omega \in \Omega$. Therefore, $(c^\pi \in \mathbb{R})_\Omega = \emptyset$. Now, note that $T_0 \vdash c \in \mathbb{R}$. Let $\delta(y) \in \mathcal{L}\{\in\}$ be a reduced defining formula for \mathbb{R} , so that $\text{ZFC}_\infty \vdash \forall y(y \in \mathbb{R} \leftrightarrow \delta(y))$. Then $\zeta \in T_0$, where $\zeta = \exists y(\delta(y) \wedge c \in y)$. We then have $\zeta^\pi = \exists y(\delta(y) \wedge c^\pi \in y)$, so that $T_0 \vdash \zeta^\pi \leftrightarrow c^\pi \in \mathbb{R}$. It follows that $\zeta_\Omega^\pi = (c^\pi \in \mathbb{R})_\Omega = \emptyset$, \mathbb{P} -a.s. Thus, $\overline{\mathbb{P}}\zeta_\Omega^\pi = 0$, contradicting Lemma 7.1.10. This shows that $c^\pi \in \{c, c_0, c_1\}$. Similar arguments show that $c_0^\pi \in \{c, c_0, c_1\}$ and $c_1^\pi \in \{c, c_0, c_1\}$. We therefore have $\pi\{c, c_0, c_1\} = \{c, c_0, c_1\}$. \square

(R:stick) Remark 7.5.5. The above proof relies on the fact that ν_0 is a continuous measure. Since we are only concerned with the relative positioning of M and c , we are free to randomize the location of M . By choosing a ν_0 which is continuous, we are randomizing M so that it has probability 0 of sitting at any fixed location. This allows us to narrow down the possible permutations in the principle of indifference, and thereby simplify our proofs. Without this construction, the results still hold, but the proofs would be more complicated. Namely, we would need to consider the possibility that M and B are interchanged for some interval B .

In Section 7.2.2, we saw how to deal with this extra complication in the finite setting. There, we needed to deal with the possibility that the symbol for the result of the coin toss, c , was interchanged with one of the symbols for heads and tails.

(L:stick-3) Lemma 7.5.6. *Let $P(\varphi \mid X) = p$ and $X^\pi \in \text{ante } P$. Assume π is not the identity. Then $\mathbf{s}^\pi = \mathbf{s}$ for all $\mathbf{s} \in L$, except for the following:*

- (a) $c_0^\pi = c_1$,
- (b) $c_1^\pi = c_0$, and
- (c) $B_*^\pi = (\rho B)_*$ for all $B \in \mathcal{E}$.

Proof. By Lemma 7.5.4, we have $\mathbf{s}^\pi = \mathbf{s}$ for all $\mathbf{s} \in L_{\text{ZFC}}$. We first show that $M^\pi = M$. Let $\mathbf{s} = M^\pi$ and assume $\mathbf{s} \neq M$. By Lemma 7.5.4, we have $\mathbf{s} = B_*$

for some $B \in \mathcal{E}$. Let $\zeta = B_* \subseteq M \wedge B_* \neq M \in T_0$. By Lemma 7.5.4, we must have $M^{-\pi} = B'_*$ for some $B' \in \mathcal{E}$. Then $\zeta^{-\pi} = M \subseteq B'_* \wedge M \neq B'_*$. But $\omega \models \neg M \subseteq B'_*$ for all $\omega \in \Omega$. Therefore, $\zeta_\Omega^{-\pi} = \emptyset$, contradicting Lemma 7.1.10. Hence, $M^\pi = M$.

We next show that $c^\pi = c$. Let $\mathfrak{s} = c^\pi$ and assume $\mathfrak{s} \neq c$. By Lemma 7.5.4, we have $\mathfrak{s} \in \{c_0, c_1\}$. Suppose $\mathfrak{s} = c_0$. Let

$$\varphi(x, y) = (\exists uv \in \mathbb{R})(u \neq \underline{0} \wedge y = f_{uv}^{\text{Img}}([0, 1]) \wedge (x = u \vee x = u + v)).$$

Then $\varphi(x, y)$ says that x is an endpoint of y . Therefore, $\zeta = \varphi^{\text{rd}}(c_0, M) \in T_0$. Since $\zeta^{-\pi} = \varphi^{\text{rd}}(c, M)$, we have

$$\begin{aligned} h^{-1}\zeta_\Omega^{-\pi} &= \{(r, t, n) \in S \mid r = 0 \text{ or } r = 1\} \\ &= \{0\} \times \mathbb{R} \times \{0, 1\} \cup \{1\} \times \mathbb{R} \times \{0, 1\}. \end{aligned}$$

Since ν_0 is continuous, this gives $\overline{\mathbb{P}}\zeta_\Omega^{-\pi} = \overline{v}h^{-1}\zeta_\Omega^{-\pi} = 0$, contradicting Lemma 7.1.10. Hence, $\mathfrak{s} \neq c_0$. A similar argument shows $\mathfrak{s} \neq c_1$. Thus, $c^\pi = c$.

We next show that $c_0^\pi = c_1$. Assume not. Then, by Lemma 7.5.4, we have $c_0^\pi = c_0$ and $c_1^\pi = c_1$. We will show that $B_*^\pi = B_*$ for all $B \in \mathcal{E}$, contradicting the assumption that π is not the identity permutation. Let $B \in \mathcal{E}$ and let $\mathfrak{s} = B_*^\pi$. Since $M^\pi = M$, Lemma 7.5.4 implies $\mathfrak{s} = B'_*$ for some $B' \in \mathcal{E}$. Note that $\varphi_B^{\text{rd}}(x) \in \mathcal{L}\{\in, M, c_0\}$. Let $\zeta = \varphi_B^{\text{rd}}(B_*) \in T_0$. Since $M^\pi = M$ and $c_0^\pi = c_0$, we have $\zeta^\pi = \varphi_B^{\text{rd}}(B'_*)$. Lemma 7.1.10 implies $\zeta_\Omega^\pi \neq \emptyset$. Hence, we may choose $x = (r, t, n) \in S$ such that $\omega \models \zeta^\pi$, where $\omega = \omega^x$. Therefore, there exists a and b in the domain of ω such that $\omega \models \varphi_M[a, b]$, $c_0^\omega = b$, and $\omega \models (B'_* = f_{uv}^{\text{Img}}(\underline{B}))[a, b]$. By the construction of ω , we have

$$\omega \models (\forall uv \in \mathbb{R})(\varphi_M(u, v) \rightarrow u = \underline{1} \vee u = \underline{-1}).$$

Hence, $a \in \{\underline{1}^{\omega_0}, \underline{-1}^{\omega_0}\}$.

Suppose $n = 0$. Then $b = c_0^\omega = \underline{t}^{\omega_0}$. If $a = \underline{-1}^{\omega_0}$, then

$$\omega \models (M = f_{uv}^{\text{Img}}([0, 1]))[\underline{-1}^{\omega_0}, \underline{t}^{\omega_0}],$$

which implies $M^\omega = [\underline{t-1}, \underline{t}]^{\omega_0}$. But $M^\omega = [\underline{t}, \underline{t+1}]^{\omega_0}$, and so we have $a = \underline{1}^{\omega_0}$. Thus, $\omega \models (B'_* = f_{uv}^{\text{Img}}(\underline{B}))[\underline{1}^{\omega_0}, \underline{t}^{\omega_0}]$. But $\omega_0 \models (\tau_t = f_{uv})[\underline{1}^{\omega_0}, \underline{t}^{\omega_0}]$ and also $\omega_0 \models \tau_t^{\text{Img}}(\underline{B}) = \tau_t B$. Hence, $\omega \models B'_* = \tau_t \underline{B}$, which gives $(B'_*)^\omega = \tau_t \underline{B}^{\omega_0}$. On the other hand, by the definition of ω , and since $n = 0$, we have $(B'_*)^\omega = (\tau_t B')^{\omega_0}$. Therefore, $(\tau_t B')^{\omega_0} = \tau_t \underline{B}^{\omega_0}$, which implies $\tau_t B' = \tau_t B$, so that $B' = B$. Hence, $\mathfrak{s} = B'_* = B_*$. A similar proof in the case $n = 1$ also yields $\mathfrak{s} = B_*$. Thus, $B_*^\pi = B_*$, completing the proof that $c_0^\pi = c_1$.

By Lemma 7.5.4, we must have $c_1^\pi = c_0$, so that both (a) and (b) hold. For (c), let $B \in \mathcal{E}$ and let $\mathfrak{s} = B_*^\pi$. As above, $\mathfrak{s} = B'_*$ for some $B' \in \mathcal{E}$. Also as above, let $\zeta = \varphi_B^{\text{rd}}(B_*) \in T_0$. Then $\omega \models \zeta^\pi$ if and only if $\omega \models \varphi_B(B_*)^\pi$, and

$$\varphi_B(B_*)^\pi = (\exists uv \in \mathbb{R})(\varphi_M(u, v) \wedge c_1 = v \wedge B'_* = f_{uv}^{\text{Img}}(\underline{B})).$$

The above argument, with c_1 in place of c_0 , shows that $B' = \rho B$. \square

Proof of main result

With all of the above preparation, we are now ready to prove Theorem 7.5.2.

Proof of Theorem 7.5.2. Let P be the inductive theory defined above. Since

$$(c \in B_*)_{\Omega} = B \times \mathbb{R} \times \{0\} \cup \rho B \times \mathbb{R} \times \{1\},$$

we have $P(c \in B_* \mid T_0) = \overline{\mathbb{P}}(c \in B_*)_{\Omega} = (\nu_0 B + \nu_0 \rho B)/2$. Hence, P satisfies (i) in the definition of \mathcal{C} . Suppose $P(\varphi \mid X) = p$ and $X^{\pi} \in \text{ante } P$. If π is the identity, then it is trivially the case that $P(\varphi^{\pi} \mid X^{\pi}) = p$. Assume π is not the identity. Define $g : S \rightarrow S$ by $g(r, t, n) = (r, t, 1 - n)$. Then g is a pointwise isomorphism from (S, Γ, ν) to itself. Recall that h is the function that maps $x = (r, t, n)$ to ω^x . Recall also the function h_{π} , used in Section 7.1.4 in the construction of \mathcal{P}^{π} . By Lemma 7.5.6, we have $h_{\pi} = h \circ g \circ h^{-1}$. Thus, $\Omega^{\pi} = \Omega$ and

$$\mathbb{P} \circ h_{\pi}^{-1} = \mathbb{P} \circ h \circ g \circ h^{-1} = \nu \circ g \circ h^{-1} = \nu \circ h^{-1} = \mathbb{P},$$

so that $\mathcal{P}^{\pi} = \mathcal{P}$. By Theorem 7.1.7, we have $\mathcal{P} \models (X^{\pi}, \varphi^{\pi}, p)$, so that $P(\varphi^{\pi} \mid X^{\pi}) = p$. This shows that P satisfies (ii), and hence, $P \in \mathcal{C}$. Therefore, \mathcal{C} is consistent.

Now let $P \in \mathcal{C}$. Let π be the permutation described in Lemma 7.5.6. Since π fixes L_{ZFC} , we have $\text{ZFC}^{\pi} = \text{ZFC}$. Since $M^{\pi} = M$ and $c^{\pi} = c$, we have $\varphi_1^{\pi} = \varphi_1$ and $\varphi_2^{\pi} = \varphi_2$. Now,

$$\begin{aligned} \text{ZFC}, u \in \mathbb{R}, v \in \mathbb{R} \vdash (\varphi_M(u, v) \wedge c_0 = v \wedge c_1 = u + v) \\ \leftrightarrow (\varphi_M(-u, u + v) \wedge c_1 = u + v \wedge c_0 = -u + (u + v)). \end{aligned}$$

Therefore, $\text{ZFC}_{\infty} \vdash \varphi_3^{\pi} \leftrightarrow \varphi_3$. Similarly,

$$\begin{aligned} \text{ZFC}, u \in \mathbb{R}, v \in \mathbb{R}, \varphi_3 \vdash (\varphi_M(u, v) \wedge c_0 = v \wedge B_* = f_{uv}^{\text{Img}}(\underline{B})) \\ \leftrightarrow (\varphi_M(-u, u + v) \wedge c_1 = u + v \wedge (\rho B)_* = f_{-u, u+v}^{\text{Img}}(\underline{B})). \end{aligned}$$

Therefore, $\text{ZFC}, \varphi_3 \vdash \varphi_B(B_*)^{\pi} \leftrightarrow \varphi_B(B_*)$. Altogether, this implies $T_0^{\pi} = T_0$, so that

$$P(c \in B_* \mid T_0) = P((c \in B_*)^{\pi} \mid T_0) = P(c \in (\rho B)_* \mid T_0),$$

by the principle of indifference. \square

7.5.3 Adding a defined constant

(S: [0, 1]-defn) Let us return to example of Section 7.5.1, in which we have a constant c , about which we only know that $c \in [0, 1]$. We saw that in this case, the principle of indifference has nothing to say about the distribution of c .

From here, let us proceed as in Section 7.3.4. That is, let us expand our language by defining $d = 1 - c$. After making this seemingly harmless addition, we turn our attention back to c , and ask again what the principle of indifference

has to say. This time, we find the same result that we obtained in (7.5.1). That is, the distribution of c must be symmetric under the reflection $r \mapsto 1 - r$.

However, as in Section 7.3.4, this result is misleading. In adding d to our language, we are also compelled to expand T_0 to a larger theory T'_0 which includes the definition of d . In the expanded theory T'_0 , the implicit meaning of c has changed. In T'_0 , we can no longer determine which of c and d is the original number, and which is its reflection. Just as in Section 7.3.4, it is as if c is equally likely to be the original as it is to be the reflection. It is no surprise, then, that c must have a symmetric distribution.

In other words, by adding d , we have changed the problem. We no longer have a single unknown number in $[0, 1]$. We now have two unknown numbers that are reflections of one another, and we are unable to definitively identify the original. We have simultaneously added information (by adding a second number) and lost information (by losing track of which was the original). So although (7.5.1) still holds, the situation is not the same. In this case, (7.5.1) is answering a different question.

To make this precise, let c and d be constant symbols not in L_{ZFC} and let $L = L_{\text{ZFC}} \cup \{c, d\}$. Let $f(r) = 1 - r$ and define

$$T_0 = \text{ZFC} + \{c \in [0, 1], d = f(c)\}.$$

Note that we did not bother to express things in terms of a definitorial extension. It is the case, however, that T_0 is a definitorial extension of $\text{ZFC} + c \in [0, 1]$, and d is defined by $\forall y(y = d \leftrightarrow y = f(c))$.

Let \mathcal{C} be the inductive condition consisting of all inductive theories $P \subseteq \mathcal{L}^{\text{IS}}$ with root T_0 such that

- (i) c is Borel given T_0 , and
- (ii) P satisfies the principle of indifference.

Proposition 7.5.7. *The condition \mathcal{C} is consistent. Moreover, if $P \in \mathcal{C}$, then*

$$P(c \in \underline{V} \mid T_0) = P(c \in \underline{f^{\text{Img}}(\underline{V})} \mid T_0), \quad (7.5.4) \quad \boxed{\text{stick-defn}}$$

for all $V \in \mathcal{B}([0, 1])$.

Proof. Let $S = [0, 1]$ and $\Gamma = \mathcal{B}([0, 1])$. Let ν be a probability measure on (S, Γ) such that $\nu \rho V = \nu V$ for all $V \in \Gamma$, and ν is continuous. That is, $\nu\{r\} = 0$ for all $r \in S$. Let ω_0 be an L_{ZFC} -structure such that $\omega_0 \models \text{ZFC}$. For each $r \in S$, let $\omega = \omega^r$ be the L -expansion of ω_0 given by $c^\omega = \underline{r^{\omega_0}}$ and $d^\omega = \underline{1 - r^{\omega_0}}$. Let $\Omega = \{\omega^r \mid r \in S\}$, let $h : S \rightarrow \Omega$ denote the function $r \mapsto \omega^r$, and let $\mathcal{P} = (\Omega, \Sigma, \mathbb{P})$ be the measure space image of (S, Γ, ν) under h . Then $\mathcal{P} \models T_0$, so we may define the inductive theory $P = \mathbf{Th} \mathcal{P} \upharpoonright_{[T_0, \text{Th} \mathcal{P}]}$.

By construction, c is Borel given T_0 . Suppose $P(\varphi \mid X) = p$ and $X^\pi \in \text{ante } P$. If π is the identity, then it is trivially the case that $P(\varphi^\pi \mid X^\pi) = p$. Assume, then, that π is not the identity.

Let $\mathbf{s} = c^\pi$. Assume $\mathbf{s} \in L_{\text{ZFC}}$. By the argument following (7.4.1), we have $P(\mathbf{s} = c \mid X) = 1$. Note that $h^{-1}(\mathbf{s} = c)_\Omega = \{r \in S \mid \mathbf{s}^{\omega_0} = \underline{r^{\omega_0}}\}$.

Since $r \neq r'$ implies $\text{ZFC} \vdash \underline{r} \neq \underline{r}'$, and since $\omega_0 \equiv \text{ZFC}$, there can be at most one $r \in \mathbb{R}$ such that $\mathfrak{s}^{\omega_0} = \underline{r}^{\omega_0}$. Since ν is continuous, this implies $P(\mathfrak{s} = c \mid T_0) = \bar{\nu} h^{-1}(\mathfrak{s} = c)_\Omega = 0$, which contradicts $P(\mathfrak{s} = c \mid X) = 1$. This shows that $\mathfrak{s} \notin L_{\text{ZFC}}$. Similarly, $d^\pi \notin L_{\text{ZFC}}$. It follows from (7.4.1) that $\mathfrak{s}^\pi = \mathfrak{s}$ for all $\mathfrak{s} \in L^\infty$. Also, since π is not the identity, we have $c^\pi = d$ and $d^\pi = c$.

Now define $g : S \rightarrow S$ by $gr = 1 - r$. Then g is a pointwise isomorphism from (S, Γ, ν) to itself and $h_\pi = h \circ g \circ h^{-1}$. As in the proof of Theorem 7.5.2, it follows that $\mathcal{P}^\pi = \mathcal{P}$. Hence, by Theorem 7.1.7, we have $\mathcal{P} \models (X^\pi, \varphi^\pi, p)$, so that $P(\varphi^\pi \mid X^\pi) = p$. This shows that P satisfies the principle of indifference, and hence, $P \in \mathcal{C}$. Therefore, \mathcal{C} is consistent.

Now let $P \in \mathcal{C}$. Let π be the permutation described above. Then $T_0^\pi = T_0$, so that

$$P(c \in \underline{V} \mid T_0) = P((c \in \underline{V})^\pi \mid T_0) = P(d \in \underline{V} \mid T_0),$$

by the principle of indifference. But $T_0 \vdash d \in \underline{V} \leftrightarrow c \in \underline{f}^{\text{Img}}(\underline{V})$. Hence, (7.5.4) follows from the rule of logical implication and Proposition 3.2.14. \square

There is nothing special about the function $f(r) = 1 - r$ in this example. What is essential is that f is measurable and $f \circ f = \iota$, where ι is the identity. For example, let

$$f(x) = \begin{cases} 1 - 2x & \text{if } 0 \leq x \leq \frac{1}{3}, \\ \frac{1}{2} - \frac{1}{2}x & \text{if } \frac{1}{3} < x \leq 1. \end{cases}$$

If we change the definition of d from $d = 1 - c$ to $d = \underline{f}(c)$, then we can adapt the above proof so that we obtain (7.5.4) in this case as well. However, the resulting distribution of c would no longer be symmetric under reflection. Instead, it would be symmetric under f . In particular, after defining $d = \underline{f}(c)$, the principle of indifference would tell us that

$$P(c \in [0, 1/3] \mid T_0) = P(c \in [1/3, 1] \mid T_0) = 1/2.$$

This is decidedly different from the situation obtained when $f(r) = 1 - r$. We can therefore see clearly that the act of defining d is not a harmless one. As described above and in Section 7.3.4, when we introduce d via definition, we change the background assumptions in T_0 , which in turn changes the meaning of c . The principle of indifference, therefore, can produce different distributions for c , depending on the definition of d .

7.6 Examples in the plane

(S:indiff-plane) In this section, we present examples of the principle of indifference that are situated in the Euclidean plane, \mathbb{R}^2 .

7.6.1 A point on a circle

(S:pt-on-circ) In our first example in this section, we consider a circle of some unspecified diameter, and we let c be a point that lies somewhere along the circumference of the circle. As with a point on a rod in Section 7.5.2, we do not simply want to replace the circle in the statement of the problem with the unit circle, $\mathbb{S}^1 = \{(x, y) \in \mathbb{R}^2 \mid x^2 + y^2 = 1\}$. Instead, we want to replace the circle with a subset $M \subseteq \mathbb{R}^2$ that can be parameterized in the usual way by the radian angles in $[0, 2\pi)$.

To make this precise, define $e : \mathbb{R} \times \mathbb{R}^2 \rightarrow \mathbb{R}^2$ by

$$e(\theta, r) = \begin{pmatrix} \cos \theta & -\sin \theta \\ \sin \theta & \cos \theta \end{pmatrix} \begin{pmatrix} r_1 \\ r_2 \end{pmatrix},$$

so that e rotates r counterclockwise by an angle of θ radians. For $t \in \mathbb{R}^2$ and $B \subseteq \mathbb{R}^2$, let $t + B = \{t + r \mid r \in B\}$. For $\theta \in \mathbb{R}$, let $[\theta] = \theta - 2\pi \lfloor \theta/2\pi \rfloor$, where $\lfloor \theta/2\pi \rfloor$ is the greatest integer less than or equal to $\theta/2\pi$. Then $[\theta] \in [0, 2\pi)$ and $e([\theta], r) = e(\theta, r)$ for all $r \in \mathbb{R}^2$.

In ZFC, we use $\underline{\mathbb{R}}^2$ as shorthand for $\mathbb{R} \times \mathbb{R}$. We extend $+$ in ZFC so that it also denotes vector addition in $\underline{\mathbb{R}}^2$, and we extend \cdot so that it also denotes scalar multiplication. As we did for \mathbb{R} and $\mathcal{B}(\mathbb{R})$, we add an explicitly defined symbol \underline{r} for each $r \in \mathbb{R}^2$, and an explicitly defined symbol \underline{V} for each $V \in \mathcal{B}(\mathbb{R}^2)$. If $h : \mathbb{R}^m \rightarrow \mathbb{R}^n$ is Borel measurable, where $m, n \in \{1, 2\}$, then we can explicitly define \underline{h} in ZFC.

We now create our extralogical signature L and our root T_0 as we did in Section 7.5.2. Define $\delta(u, v, w, y) \in \mathcal{L}_{\text{ZFC}}$ by

$$\begin{aligned} \delta(u, v, w, y) = & (u \notin \underline{\mathbb{R}} \vee v \notin \underline{\mathbb{R}}^2 \vee w \notin \underline{\mathbb{R}}) \wedge y = \emptyset \\ & \vee u \in \underline{\mathbb{R}} \wedge v \in \underline{\mathbb{R}}^2 \wedge w \in \underline{\mathbb{R}} \wedge y \in (\underline{\mathbb{R}}^2)^{\mathbb{R}^2} \\ & \wedge (\forall x \in \underline{\mathbb{R}})(y(x) = u \cdot \underline{e}(w, x) + v). \end{aligned}$$

Then $\text{ZFC} \vdash \forall uvw \exists! y \delta(u, v, w, y)$. Hence, we could explicitly define the function symbol F by $y = Fuvw \leftrightarrow \delta(u, v, w, y)$, and then let f_{uvw} be shorthand for the term $Fuvw$. We do not, in fact, add the symbol F to our extralogical signature, but instead regard both F and f_{uvw} as shorthand. We also adopt the shorthand, $\text{dom}(f)$ and $f^{\text{Img}}(z)$, given by

$$\begin{aligned} \text{dom}(f) &= \{x \in \underline{\mathbb{R}}^2 \mid (\exists y \in \underline{\mathbb{R}}^2)(x, y) \in f\}, \\ f^{\text{Img}}(z) &= \{f(x) \mid x \in z \cap \text{dom}(f)\}. \end{aligned}$$

Let $\mathcal{E} = \{B \in \mathcal{B}(\mathbb{S}^1) \mid \emptyset \subset B \subset \mathbb{S}^1\}$ and let

$$C = \{M, c\} \cup \{B_* \mid B \in \mathcal{E}\}$$

be a set of distinct constant symbols not in L_{ZFC} . Let $L = L_{\text{ZFC}}C$.

Let $\varphi_c = c \in M$ and

$$\varphi_M(u, v, w) = u > \underline{0} \wedge w \in \underline{[0, 2\pi)} \wedge M = f_{uvw}^{\text{Img}}(\underline{\mathbb{S}^1}).$$

For $B \in \mathcal{E}$, define

$$\begin{aligned} \varphi_B &= (\exists uw \in \mathbb{R})(\exists v \in \mathbb{R}^2) \\ &(\varphi_M(u, v, w) \wedge \{(1, 0)\}_* = f_{uvw}^{\text{Img}}(\{(1, 0)\}) \wedge B_* = f_{uvw}^{\text{Img}}(\underline{B})), \end{aligned}$$

and let $T_0 = \text{ZFC} + \{\varphi_c\} \cup \{\varphi_B \mid B \in \mathcal{E}\}$.

In this presentation, we have streamlined our construction of T_0 , compared to what was done in the rod example of Section 7.5.2. For instance, we do not have a separate sentence which says that M is a circle. Rather, that fact is contained in each sentence φ_B . The sentence φ_c says that c is a point on the circle M . And the sentences φ_B say that B_* is a Borel subset of M that is geometrically related to $\{(1, 0)\}_*$ in the same way that B is related to $\{(1, 0)\}$. In the rod example, we first named our endpoints and then named our Borel sets relative to them. Here, we are first naming a point to serve as the initial and terminal point of a parameterization, and then naming our Borel sets relative to that point.

Now let \mathcal{C} be the inductive condition consisting of all inductive theories $P \subseteq \mathcal{L}^{\text{IS}}$ with root T_0 such that

- (i) $P(c \in B_* \mid T_0)$ exists for all $B \in \mathcal{E}$, and
- (ii) P satisfies the principle of indifference.

(P:circle) **Proposition 7.6.1.** *The condition \mathcal{C} is consistent and, for every $B \in \mathcal{E}$, we have $\mathbf{P}_c(c \in B_* \mid T_0) = m(B)$, where m is the uniform measure on \mathbb{S}^1 .*

The proof of Proposition 7.6.1 will come at the end of this subsection. The proof follows the same lines as the proof of Theorem 7.5.2. We first prove consistency by building an \mathcal{L} -model, $\mathcal{P} = (\Omega, \Sigma, \mathbb{P})$, as follows. Let ν_0 be a probability measure on $(\mathbb{S}^1, \mathcal{B}(\mathbb{S}^1))$ such that ν_0 is continuous. That is, $\nu_0\{r\} = 0$ for all $r \in \mathbb{S}^1$. Let $S = \mathbb{S}^1 \times \mathbb{S}^1 \times [0, 2\pi)$, $\Gamma = \mathcal{B}(S)$, and $\nu = \nu_0 \times \nu_0 \times m_0$, where m_0 is the uniform measure on $[0, 2\pi)$.

Let ω_0 be an L_{ZFC} -structure such that $\omega_0 \models \text{ZFC}$. If $x = (r, t, \theta) \in S$, then let $\omega = \omega^x$ be the L -expansion of ω_0 defined by $M^\omega = \underline{t + \mathbb{S}^{1\omega_0}}$, $c^\omega = \underline{t + r^{\omega_0}}$, and $B_*^\omega = \underline{t + e(\theta, B)^{\omega_0}}$.

Let $\Omega = \{\omega^x \mid x \in S\}$, let $h : S \rightarrow \Omega$ denote the function $x \mapsto \omega^x$, and let \mathcal{P} be the measure space image of (S, Γ, ν) under h .

By a proof similar to that of Lemma 7.5.3, we have $\mathcal{P} \models T_0$. We may therefore define the inductive theory $P = \mathbf{Th} \mathcal{P} \upharpoonright_{[T_0, \text{Th} \mathcal{P}]}$.

(L:circle) **Lemma 7.6.2.** *Let $P(\varphi \mid X) = p$ and $X^\pi \in \text{ante } P$. Then $\mathbf{s}^\pi = \mathbf{s}$ for all $\mathbf{s} \in L_{\text{ZFC}} \cup \{M, c\}$ and there exists $\theta_0 \in [0, 2\pi)$ such that $B_*^\pi = e(\theta_0, B)_*$ for all $B \in \mathcal{E}$.*

Proof. Applying the methods used in the proof of Lemma 7.5.4 and the first part of the proof of Lemma 7.5.6, we obtain that $\mathbf{s}^\pi = \mathbf{s}$ for all $\mathbf{s} \in L_{\text{ZFC}} \cup \{M, c\}$. Hence, we must have $\{(1, 0)\}_*^\pi = B_*^0$ for some $B^0 \in \mathcal{E}$. Let $\zeta = (\exists! y \in \mathbb{R}^2)(y \in \{(1, 0)\}_*) \in T_0$. By Lemma 7.1.10, the set ζ_Ω^π is nonempty. Choose $\omega \in \zeta_\Omega^\pi$.

Then $\omega \equiv (\exists! y \in \mathbb{R}^2)(y \in B_*^0)$. Write $\omega = \omega^x$, where $x = (r, t, \theta) \in S$. Then there exists a unique b in the domain of ω_0 such that $b \in^{\omega_0} (\mathbb{R}^2)^{\omega_0}$ and $b \in^{\omega_0} (B_*^0)^\omega = \underline{t + e(\theta, B^0)^{\omega_0}}$. Thus, $\omega_0 \equiv (\exists! y \in \mathbb{R}^2)(y \in \underline{t + e(\theta, B^0)})$. This implies $|B^0| = \bar{1}$, so that we may write $B^0 = \{r_0\}$ for some $r_0 \in \mathbb{S}^1$. We therefore have $\{(1, 0)\}_*^\pi = \{r_0\}_*$.

Now let $B \in \mathcal{E}$ be arbitrary. Then $B_*^\pi = B'_*$ for some $B' \in \mathcal{E}$. By Lemma 7.1.10, we may choose $x = (r, t, \theta) \in S$ such that $\omega = \omega^x \equiv \varphi_B^\pi$. Note that

$$\begin{aligned} \varphi_B^\pi &\equiv_{\text{ZFC}} (\exists! uv \in \mathbb{R})(\exists! v \in \mathbb{R}^2) \\ &\quad (\varphi_M(u, v, w) \wedge \{r_0\}_* = f_{uvw}^{\text{Img}}(\{(1, 0)\}) \wedge B'_* = f_{uv}^{\text{Img}}(\underline{B})). \end{aligned}$$

Choose θ_0 such that $r_0 = e(\theta_0, (1, 0))$. Then

$$\{r_0\}_*^\omega = \underline{t + e(\theta, \{r_0\})}^{\omega_0} = \underline{t + e([\theta_0 + \theta], \{(1, 0)\})}^{\omega_0},$$

so that

$$\omega \equiv (\{r_0\}_* = f_{uvw}^{\text{Img}}(\{(1, 0)\}))[\underline{1}^{\omega_0}, \underline{t}^{\omega_0}, [\underline{\theta_0 + \theta}]^{\omega_0}].$$

Since $\omega \equiv \varphi_B^\pi$, we must have

$$\omega \equiv (B'_* = f_{uv}^{\text{Img}}(\underline{B}))[\underline{1}^{\omega_0}, \underline{t}^{\omega_0}, [\underline{\theta_0 + \theta}]^{\omega_0}],$$

which implies

$$\underline{t + e(\theta, B')^{\omega_0}} = \underline{t + e([\theta_0 + \theta], B)^{\omega_0}},$$

and therefore $B' = e(\theta_0, B)$. \square

Proof of Proposition 7.6.1. Let P be the inductive theory defined above. If $x = (r, t, \theta)$, then

$$\begin{aligned} \omega^x \equiv c \in B_* &\quad \text{iff } \underline{t + r}^{\omega_0} \in^{\omega_0} \underline{t + e(\theta, B)^{\omega_0}} \\ &\quad \text{iff } \omega_0 \equiv \underline{t + r} \in \underline{t + e(\theta, B)} \\ &\quad \text{iff } t + r \in t + e(\theta, B) \\ &\quad \text{iff } r \in e(\theta, B). \end{aligned}$$

Thus, $h^{-1}(c \in B_*)_\Omega = \{(r, t, \theta) \mid r \in e(\theta, B)\}$. Note that $r \in e(\theta, B)$ if and only if $e(-\theta, r) \in B$. Hence,

$$\begin{aligned} \nu h^{-1}(c \in B_*)_\Omega &= \int_S 1_{h^{-1}(c \in B_*)_\Omega} d\nu \\ &= \int_{\mathbb{S}^1} \int_0^{2\pi} 1_B(e(-\theta, r)) m_0(d\theta) \nu_0(dr) \\ &= \int_{\mathbb{S}^1} m(B) \nu_0(dr) = m(B). \end{aligned}$$

Since $\mathbb{P} = \nu \circ h^{-1}$, it follows that P satisfies (i) in the definition of \mathcal{C} .

Suppose $P(\varphi \mid X) = p$ and $X^\pi \in \text{ante } P$. Let θ_0 be as in Lemma 7.6.2 and define $g : S \rightarrow S$ by $g(r, t, \theta) = (r, t, [\theta_0 + \theta])$. We may use g as in the proof of Theorem 7.5.2 to show that $\mathcal{P}^\pi = \mathcal{P}$, so that Theorem 7.1.7 yields $P(\varphi^\pi \mid X^\pi) = p$. This shows that P satisfies (ii), and hence, $P \in \mathcal{C}$. Therefore, \mathcal{C} is consistent.

Now let $P \in \mathcal{C}$ be arbitrary. Given $\theta_0 \in \mathbb{R}$, let π be the L -permutation satisfying $\mathfrak{s}^\pi = \mathfrak{s}$ for all $\mathfrak{s} \in L_{\text{ZFC}} \cup \{M, c\}$ and $B_*^\pi = e(\theta_0, B)_*$ for all $B \in \mathcal{E}$. Then $T_0^\pi = T_0$, so by the principle of indifference, we have

$$P(c \in B_* \mid T_0) = P(c \in e(\theta_0, B)_* \mid T_0) \quad (7.6.1) \quad \boxed{\text{circle}}$$

for all $B \in \mathcal{E}$. Let us adopt the shorthand notation, $\emptyset_* = \emptyset$ and $\mathbb{S}_*^1 = M$. Since $P(c \in B_* \mid T_0)$ exists for all $B \in \mathcal{E}$, we may define $m' : \mathcal{B}(\mathbb{S}^1) \rightarrow [0, 1]$ by $m'(B) = P(c \in B_* \mid T_0)$. Note that if $B, B' \in \mathcal{B}(\mathbb{S}^1)$ and $B \cap B' = \emptyset$, then $T_0 \vdash \neg(c \in B_* \wedge c \in B'_*)$. Hence, Theorem 3.2.24 implies that m' is a probability measure on $(\mathbb{S}^1, \mathcal{B}(\mathbb{S}^1))$. By (7.6.1), we have $m'(B) = m'(e(\theta_0, B))$ for all $B \in \mathcal{B}(\mathbb{S}^1)$ and all $\theta_0 \in \mathbb{R}$. This implies $m' = m$, so that $P(c \in B_* \mid T_0) = m(B)$. Since P was arbitrary, $\mathbf{P}_{\mathcal{C}}(c \in B_* \mid T_0) = m(B)$. \square

7.6.2 Bertrand's paradox

Introduction

In 1888, Joseph Bertrand posed the following problem (see [2]). Consider an equilateral triangle inscribed in a circle. Let ℓ be a chord of the circle, chosen at random. What is the probability that the chord is longer than a side of the triangle?

It is considered a “paradox” because Bertrand presented three different solutions, all purporting to use the principle of indifference, that gave three different answers: $1/3$, $1/2$, and $1/4$. Of course, this is only “paradoxical” if we have the prior expectation that the principle of indifference ought to produce a unique answer. We have already seen, however, that this is not always the case. The principle of indifference is a tool that can narrow down the possible distributions in certain circumstances, but it does not necessarily determine for us a unique distribution. In the example of the rod from Section 7.5.2, for instance, the principle tells us that the distribution must be symmetric under reflection. But beyond that, it leaves open a whole range of possibilities.

Something similar happens with Bertrand's chord. We will show below that, according to the principle of indifference, the distribution of the chord must be rotationally invariant. But beyond that, it has nothing more to say. Hence, all three of Bertrand's solutions (which are all rotationally invariant) are consistent with the principle of indifference. But so are many distributions that Bertrand did not consider. In fact, in Theorem 7.6.3, we show that for any $p \in [0, 1]$, it is consistent with the principle of indifference to say that the answer is p .

Notation in ZFC

To precisely formulate the problem, we first establish some new notation and shorthand. Let $\mathbb{D} = \{(x, y) \in \mathbb{R}^2 \mid x^2 + y^2 \leq 1\}$. Using the notation of Section 7.6.1, define

$$\varphi_{\mathbb{D}}^{\text{par}}(x, u, v, w) : u > 0 \wedge w \in [0, 2\pi) \wedge x = f_{uvw}^{\text{img}}(\mathbb{D})$$

Then $\varphi_{\mathbb{D}}^{\text{par}}(x, u, v, w)$ says that x is a disk in the plane and that u, v , and w are the constants in a parameterization of x . Also define

$$\varphi_{\mathbb{D}}(x) : (\exists uvw \in \mathbb{R})(\exists v \in \mathbb{R}^2)\varphi_{\mathbb{D}}^{\text{par}}(x, u, v, w)$$

Then $\varphi_{\mathbb{D}}(x)$ simply says that x is a disk in the plane.

Let $\varphi_{\text{tri}}(x)$ be a formula which says that x is a nondegenerate equilateral triangle in the plane. The exact details of this formula are not relevant for our purposes and will be omitted. In what follows, we will similarly omit the details of other formulas whose descriptions are given only verbally.

Let $\varphi_{\text{seg}}(x)$ be a formula which says that x is a line segment in the plane whose length is positive and finite. Let $\varphi_{\text{ins}}(x, y) = \varphi_{\text{tri}}(x) \wedge \varphi_{\mathbb{D}}(y) \wedge \zeta(x, y)$, where $\zeta(x, y)$ is a formula which says that the triangle x is inscribed in the circle that is the boundary of y . Similarly, let $\varphi_{\text{ch}}(x, y) = \varphi_{\text{seg}}(x) \wedge \varphi_{\mathbb{D}}(y) \wedge \zeta(x, y)$, where $\zeta(x, y)$ is a formula which says that the endpoints of the line segment x lie on the boundary of the disk y .

Define

$$\delta_{\text{len}}(x, y) : \neg\varphi_{\text{tri}}(x) \wedge \neg\varphi_{\text{seg}}(x) \wedge y = \emptyset \vee \varphi_{\text{tri}}(x) \wedge \zeta(x, y) \vee \varphi_{\text{seg}}(x) \wedge \zeta'(x, y)$$

Here, $\zeta(x, y)$ is a formula which says that y is the length of each side of the equilateral triangle x , and $\zeta'(x, y)$ is a formula which says that y is the length of the line segment x . Then $\text{ZFC} \vdash \forall x \exists! y \delta_{\text{len}}(x, y)$. We could therefore define the function symbol F by $y = Fx \leftrightarrow \delta_{\text{len}}(x, y)$ and let $\text{len}(x)$ be shorthand for the term Fx . We do not actually add F to our extralogical signature, and instead leave both F and $\text{len}(x)$ as shorthand. In this way, $\text{len}(x)$ is a function informally described by

$$\text{len}(x) = \begin{cases} \text{the length of } x & \text{if } x \text{ is a line segment,} \\ \text{the length of a side of } x & \text{if } x \text{ is an equilateral triangle,} \\ \emptyset & \text{otherwise.} \end{cases}$$

A first pass at setting up the problem

Let $C' = \{D, \tau, \ell\}$ and $L' = L_{\text{ZFC}}C'$. Define the deductive theory $T'_0 \subseteq (\mathcal{L}')^0$ by

$$T'_0 = \text{ZFC} + \varphi_{\mathbb{D}}(D) + \varphi_{\text{ins}}(\tau, D) + \varphi_{\text{ch}}(\ell, D).$$

Then T'_0 includes all facts in ZFC together with the following three assumptions: D is a disk in the plane, τ is an equilateral triangle inscribed in D , and ℓ is a chord of D .

Define the inductive condition \mathcal{C}' as the set of inductive theories $P \subseteq (\mathcal{L}')^{\text{IS}}$ with root T'_0 such that

- (i) $P(\text{len}(\ell) > \text{len}(\tau) \mid T'_0)$ exists, and
- (ii) P satisfies the principle of indifference.

Let $P \subseteq (\mathcal{L}')^{\text{IS}}$ be any inductive theory with root T'_0 . Suppose $P(\varphi \mid X) = p$ and $X^\pi \in \text{ante } P$. As in Remark 7.5.5, since we are only concerned with the relative positions of D , τ , and ℓ , we may assume that $P(D = \underline{B} \mid T'_0) = 0$ for all $B \in \mathcal{B}(\mathbb{R}^2)$. We may also make similar assumptions for τ and ℓ . It then follows as in the proof of Lemma 7.6.2 that π must be the identity permutation, so that $P(\varphi^\pi \mid X^\pi) = p$. In other words, every inductive theory in $(\mathcal{L}')^{\text{IS}}$ satisfies the principle of indifference. This means that the principle of indifference has nothing to say in this setting. It offers no restrictions, so that $P(\text{len}(\ell) > \text{len}(\tau) \mid T'_0)$ could be anything we like.

This, however, is misleading. We have omitted a critical assumption. Namely, we failed to interpret the fact that the chord is “chosen at random.” We will interpret this additional assumption as simply saying that the location of the chord has a probability distribution. We leave the exact nature of this distribution unspecified. To formulate this additional assumption, we must expand our extralogical signature and our root.

The complete setup and conclusion

Let $\mathcal{E} = \{B \in \mathcal{B}(\mathbb{D}) \mid \emptyset \subset B \subset \mathbb{D}\}$ and let

$$C = \mathcal{C}' \cup \{B_* \mid B \in \mathcal{E}\} = \{D, \tau, \ell\} \cup \{B_* \mid B \in \mathcal{E}\}.$$

Let $L = L_{\text{ZFC}}C$. For each $B \in \mathcal{E}$, define

$$\begin{aligned} \varphi_B(x) : (\exists uw \in \mathbb{R})(\exists v \in \mathbb{R}^2) \\ (\varphi_{\mathbb{D}}^{\text{par}}(D, u, v, w) \wedge \{(1, 0)\}_* = f_{uw}^{\text{img}}(\{(1, 0)\}) \wedge B_* = f_{uvw}^{\text{img}}(\underline{B})) \end{aligned}$$

Then let

$$\begin{aligned} T_0 &= T'_0 + \{\varphi_B(B_*) \mid B \in \mathcal{E}\} \\ &= \text{ZFC} + \varphi_{\mathbb{D}}(D) + \varphi_{\text{ins}}(\tau, D) + \varphi_{\text{ch}}(\ell, D) + \{\varphi_B(B_*) \mid B \in \mathcal{E}\}. \end{aligned}$$

Our root, T_0 , says that all the facts in ZFC hold. It also says that D is a disk in the plane, τ is an equilateral triangle inscribed in D , and ℓ is a chord of D . Regarding the subsets of D , it says that $\{(1, 0)\}_*$ is a singleton set on the boundary of D that serves as a fixed point of reference. The sets B_* are then Borel subsets of D that are geometrically related to $\{(1, 0)\}_*$ in the same manner as B is related to $\{(1, 0)\}$.

Define

$$\delta_{\text{mid}}(x, y) : \neg\varphi_{\text{seg}}(x) \wedge y = \emptyset \vee \varphi_{\text{seg}}(x) \wedge \zeta(x, y)$$

Here, $\zeta(x, y)$ is a formula which says that y is the midpoint of the line segment x . Then $T_0 \vdash \exists! y \delta_{\text{mid}}(\ell, y)$. We could therefore define the constant symbol c by $y = c \leftrightarrow \delta_{\text{mid}}(\ell, y)$. We do not actually add c to our extralogical signature, but instead leave it as shorthand. With this construction, c denotes the midpoint of the segment ℓ . We can therefore talk about the location of ℓ using the sentence $c \in B_*$, where $B \in \mathcal{E}$.

Let

$$B^{1/2} = \{(x, y) \in \mathbb{R}^2 \mid \sqrt{x^2 + y^2} < 1/2\} \in \mathcal{E}.$$

Then $T_0 \vdash \text{len}(\ell) > \text{len}(\tau) \leftrightarrow c \in B_*^{1/2}$. Hence, for any inductive theory $P \subseteq \mathcal{L}^{\text{IS}}$ with root T_0 , we have

$$P(\text{len}(\ell) > \text{len}(\tau) \mid T_0) = P(c \in (B^{1/2})_* \mid T_0), \quad (7.6.2) \boxed{\text{len-midp}}$$

by Proposition 3.2.14.

Define the inductive condition \mathcal{C} as the set of inductive theories $P \subseteq \mathcal{L}^{\text{IS}}$ with root T_0 such that

- (i) $P(c \in B_* \mid T_0)$ exists for all $B \in \mathcal{E}$, and
- (ii) P satisfies the principle of indifference.

Our main result now follows. For this, recall the notation $e(\theta, r)$ from Section 7.6.1.

(T:Bertrand) **Theorem 7.6.3.** *The inductive condition \mathcal{C} is consistent and every $P \in \mathcal{C}$ satisfies*

$$P(c \in B_* \mid T_0) = P(c \in e(\theta, B)_* \mid T_0) \quad (7.6.3) \boxed{\text{Bertrand}}$$

for all $B \in \mathcal{E}$ and all $\theta \in \mathbb{R}$. Moreover, for every $p \in [0, 1]$, there exists $P \in \mathcal{C}$ such that

$$P(\text{len}(\ell) > \text{len}(\tau) \mid T_0) = p.$$

Proof. We begin by proving consistency. Let $\mathbb{T} \subseteq \mathbb{R}^2$ be the equilateral triangle inscribed in \mathbb{D} with one corner situated at $(1, 0)$. Given $r \in \mathbb{D}$ with $|r| < 1$, let \mathbb{L}_r be the chord of \mathbb{D} whose midpoint is r .

Let ν_0 be a probability measure on $(\mathbb{D}, \mathcal{B}(\mathbb{D}))$ such that ν_0 is continuous. That is, $\nu_0\{r\} = 0$ for all $r \in \mathbb{D}$. Let $S = \mathbb{D} \times \mathbb{D} \times [0, 2\pi)$, $\Gamma = \mathcal{B}(S)$, and $\nu = \nu_0 \times \nu_0 \times m_0$, where m_0 is the uniform measure on $[0, 2\pi)$.

Let ω_0 be an L_{ZFC} -structure such that $\omega_0 \models \text{ZFC}$. If $x = (r, t, \theta) \in S$, then let $\omega = \omega^x$ be the L -expansion of ω_0 defined by $D^\omega = \underline{t} + \mathbb{D}^{\omega_0}$, $\tau^\omega = \underline{t} + \mathbb{T}^{\omega_0}$, $\ell^\omega = \underline{t} + \mathbb{L}_r^{\omega_0}$, and $B_*^\omega = \underline{t} + e(\theta, B)^{\omega_0}$.

Let $\Omega = \{\omega^x \mid x \in S\}$, let $h : S \rightarrow \Omega$ denote the function $x \mapsto \omega^x$, and let $\mathcal{P} = (\Omega, \Sigma, \mathbb{P})$ be the measure space image of (S, Γ, ν) under h . By a proof similar to that of Lemma 7.5.3, we have $\mathcal{P} \models T_0$. We may therefore define the inductive theory $P = \mathbf{Th} \mathcal{P} \downarrow_{[T_0, Th \mathcal{P}]}$.

For each $x = (r, t, \theta) \in S$ and $\omega = \omega^x$, it follows that $\omega \models c \in B_*$ if and only if $r \in e(\theta, B)$. Thus, $h^{-1}(c \in B_*)_\Omega = \{(r, t, \theta) \mid r \in e(\theta, B)\}$. Since e is

measurable, $P(c \in B_* \mid T_0) = \mathbb{P}(c \in B_*)_\Omega = \nu h^{-1}(c \in B_*)_\Omega$ exists, so that P satisfies (i) in the definition of \mathcal{C} .

Suppose $P(\varphi \mid X) = p$ and $X^\pi \in \text{ante } P$. Using methods like those in the proofs of Lemmas 7.5.4 and 7.5.6, it follows that $\mathbf{s}^\pi = \mathbf{s}$ for all $\mathbf{s} \in L_{\text{ZFC}} \cup \{D, \tau, \ell\}$, and there exists $\theta_0 \in [0, 2\pi)$ such that $B_*^\pi = e(\theta_0, B)_*$ for all $B \in \mathcal{E}$. Define $g : S \rightarrow S$ by $g(r, t, \theta) = (r, t, [\theta_0 + \theta])$. As in the proofs of Theorem 7.5.2 and Proposition 7.6.1, we can use g to show that $\mathcal{P}^\pi = \mathcal{P}$. It therefore follows that $P(\varphi^\pi \mid X^\pi) = p$, so that P satisfies the principle of indifference. Hence, $P \in \mathcal{C}$ and \mathcal{C} is consistent.

Now let $P \in \mathcal{C}$ be given. Fix $\theta_0 \in \mathbb{R}$ and $B \in \mathcal{E}$. Let π be the L -permutation such that $\mathbf{s}^\pi = \mathbf{s}$ for all $\mathbf{s} \in L_{\text{ZFC}} \cup \{D, \tau, \ell\}$, and $B_*^\pi = e(\theta_0, B)_*$ for all $B \in \mathcal{E}$. Then $T_0^\pi = T_0$, so (7.6.3) follows immediately from the principle of indifference.

Finally, let $p \in [0, 1]$. By (7.6.2), it suffices to show that there exists $P \in \mathcal{C}$ such that $P(c \in (B^{1/2})_* \mid T_0) = p$. Let ν_0 be a continuous probability measure on $(\mathbb{D}, \mathcal{B}(\mathbb{D}))$ such that $\nu_0(B^{1/2}) = p$. Construct P as in the first part of this proof. Since $r \in e(\theta, B)$ if and only if $e(-\theta, r) \in B$, it follows that

$$\begin{aligned} P(c \in (B^{1/2})_* \mid T_0) &= \nu h^{-1}(c \in (B^{1/2})_*)_\Omega \\ &= \int_{\mathbb{D}} \int_0^{2\pi} 1_{B^{1/2}}(e(-\theta, r)) m_0(d\theta) \nu_0(dr). \end{aligned}$$

But $1_{B^{1/2}}(e(-\theta, r)) = 1_{B^{1/2}}(r)$, so $P(c \in (B^{1/2})_* \mid T_0) = \nu_0(B^{1/2}) = p$. \square

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Index of Terms

- α -sequence, 14
- addition rule, *see* rule
- admissible, 36
- almost everywhere, 17
- almost surely, 18
- antecedent, 35
 - formula, 50
 - generalized —, 52
- assignment, 125, 126
 - image of an —, 127
- axiom, 32, 122
 - Karp's — theorem, 73, 129
- basis, 47
- Bayes' theorem, 41
- Bertrand's paradox, 214
- Boolean
 - algebra, 16
 - function, 69
 - measure space, 17
 - σ -algebra, 16
- bound renaming, 116
- C -substitution, 117
- cardinal, 15
 - limit —, 16
 - strongly inaccessible —, 16
 - successor —, 16
- cardinality, 15
- central limit theorem, 174
- closed, 45
- complete, 44
- completeness
 - deductive —, 75, 132
 - for inductive conditions, 85
 - inductive —, 83
- conditional expectation, 179
- connected, 47
 - strongly —, 46
- consequence relation, 71, 81, 85, 128
- consequent, 35
- consistent, 31, 48, 62, 121
- continuity rule, *see* rule
- contradiction, 31, 121
- countable additivity, 43
- countably axiomatizable, 46
- Dedekind cut, 161
- definitorial extension, 148, 150
- dependent, 97
- derivability relation
 - predicate —, 115
 - propositional —, 29, 48, 63
- dialog set, 96
- distribution, 171, 172, 177
- Dynkin system, *see* λ -system
- entire, 36
- equation, 108
- equivalent, 34, 59
- exchangeable, 200
- expansion, *see* structure
- expected value, 171
- explicit definition, 145, 146
- extralogical signature, 21
- extralogical symbol, 21
- falsum, 28, 121
- formula, 27, 108
 - defining —, 145, 146

- prime —, 108
- reduced —, 146
- frame of reference, 139
 - natural —, 141
 - real —, 165, 168, 169
- free eliminator, 120
- free occurrence, 114
- ground term, *see* term
- inclusion-exclusion, 40
- inconsistent, 31, 121
- independent, 97, 98
 - measure —, 99
- indicator function, 20
- induction
 - formula —, 28, 110
 - transfinite —, 14
- inductive condition, 62
 - determinate —, 63
 - indeterminate —, 63
- inductive model, *see* model
- inductive statement, 35, 121
- invariant, 182
- isomorphism theorem
 - deductive —, 127
 - inductive —, 134
- λ -system, 19
- law of large numbers, 173
- law of total probability, 180
- length, 110
- lift, 53
- measurable function, 19
- measure space, 17
 - complete —, 18
 - completion of a —, 19
 - image, 21
 - isomorphism, 20
- model, 70, 126
 - isomorphic —, 70, 127
 - strict —, 69
- multiplication rule, *see* rule
- negligible set, 18
- null set, 17
- ordinal, 13
 - arithmetic, 15
 - limit —, 14
 - successor —, 14
- π - λ theorem, 19
- π -system, 19
- Peano arithmetic, 132
- permutation
 - signature —, 182
 - variable —, 117
- power set, 15
- pre-theory, 46
- prime term, *see* term
- principle of indifference, 185
- probability, 35
 - kernel, 175
 - space, 17
- proof, 32, 122
- pushforward, 20
- rank of a formula, 110
- rank of a measurable set, 18
- reduct, *see* structure
- root, 49, 62
- rule
 - of deductive extension, 45
 - of deductive transitivity, 37
 - of inductive extension, 44
 - of logical equivalence, 36
 - of logical implication, 36
 - of material implication, 37
 - the addition —, 37, 40
 - the continuity —, 37, 42
 - the multiplication —, 37, 41
- σ -algebra, 17
 - completion of a —, 19
- σ -compactness, 30, 74, 117, 131
- satisfiable, 70, 76, 84, 126, 134
 - strictly —, 70, 125
- scope, 114
- semi-closed, 44
- sentence, 27, 111
- shrinks nicely, 175
- soundness
 - deductive —, 71, 130

- for inductive conditions, 85
- inductive —, 83
- standard structure of arithmetic, 23
- string, 24
- structure, 21
 - domain of a —, 21
 - expansion of a —, 22
 - isomorphic image of a —, 22
 - isomorphism, 22
 - reduct of a —, 22
- subformula, 28, 110
- substitution, 112
 - free —, 115
- tautology, 31, 121
- term, 107
 - Borel —, 171
 - ground —, 108
 - independent —, 172
 - integrable —, 171
 - jointly Borel —, 172
 - prime —, 107
 - real —, 171
- theory
 - deductive —, 33, 121
 - inductive —, 48
 - real inductive —, 162, 163, 169
- variable
 - bound —, 110
 - free —, 111
 - individual —, 107
 - propositional —, 27
- verum, 28, 121
- von Neumann hierarchy, 156

Index of Symbols

\mathbb{N}	13	$\text{AF}(P)$	50	$\text{var } \varphi$	110
\mathbb{N}_0	13	$Q \downarrow_{\mathcal{X}}$	51	$\text{bnd } \varphi$	110
\mathbb{N}'_0	14	$\mathcal{GA}(P_0)$	52	$\text{free } \varphi$	111
ω	14	$\mathbf{L}(P_0)$	53	\mathcal{L}^0	111
ω_1	15	$T(Q), T_Q$	62	$t(\vec{x})$	111
\mathfrak{P}	15	$\mathbf{P}(\mathcal{C}), \mathbf{P}_{\mathcal{C}}$	63	$\varphi(\vec{x})$	111
\mathbf{B}	17	$\mathcal{C} \vdash (X, \varphi, p)$	63	$\text{sym } \varphi$	111
$\sigma(\mathcal{E})$	17	$T(\mathcal{C}), T_{\mathcal{C}}$	63	$\text{con } \varphi$	111
$1_A, 1_A(\omega)$	20	\mathbf{B}^S	67	$\varphi(t/x)$	112
\mathcal{N}	22	\mathbf{B}^S	67	$X \vdash \varphi$ (in \mathcal{L})	115
\mathcal{N}'	23	\equiv	70	$X \vdash_{\text{fin}} \varphi$ (in \mathcal{L})	116
PV	27	φ_{Ω}	70	\mathcal{LC}	117
\mathcal{F}	27	$\mathscr{P} \models \varphi$	70	\top (in \mathcal{L})	121
\mathcal{F}_{fin}	27	$\text{Th } \mathscr{P}$ (in \mathcal{F})	70	\perp (in \mathcal{L})	121
\perp (in \mathcal{F})	28	$X \models \varphi$ (in \mathcal{F})	71	$T(X), T_X$ (in \mathcal{L})	121
\top (in \mathcal{F})	28	$\Delta(P, X)$	74	$\text{Taut}_{\mathcal{L}}$	121
φ^1, φ^0	28	$\mathscr{P} \models_{\mathcal{F}} (X, \varphi, p)$	76	\mathcal{L}^{IS}	121
$X \vdash \varphi$ (in \mathcal{F})	29	$\mathbf{Th } \mathscr{P}$	76	$\Lambda_{\mathcal{L}}$	122
$X \vdash_{\text{fin}} \varphi$ (in \mathcal{F})	29	$Q \models (X, \varphi, p)$	81	$\sim_{\mathcal{L}}$	122
$\text{Taut}_{\mathcal{F}}$	31	$\mathscr{P} \models \mathcal{C}$	84	$t^{\omega}[\vec{a}]$	125
$\Lambda_{\mathcal{F}}$	31	$\mathcal{C} \models (X, \varphi, p)$	85	$\omega \equiv \varphi[v]$	125
$\sim_{\mathcal{F}}$	32	$d_k(n)$	89	$\varphi[\mathbf{v}]_{\Omega}$	126
$T(X), T_X$ (in \mathcal{F})	34	$\delta(X)$	96	$\mathscr{P} \models \varphi[\mathbf{v}]$	126
$T + S$	34	\mathfrak{J}_{T_0}	103	$X \models \varphi$ (in \mathcal{L})	128
\equiv	34, 59	Var	107	PA_{fin}	133
\equiv_X	34	\mathcal{S}	107	PA_{-}	133
$[T_0, T_1]$	35	\mathcal{T}	107	PA	133
\mathcal{F}^{IS}	35	$\text{var } t$	107	$\text{Th } \mathscr{P}$ (in \mathcal{L})	134
(X, φ, p)	35	$\text{sym } t$	108	$\mathscr{P} \models_{\mathcal{L}} (X, \varphi, p)$	134
$P(\varphi \mid X) = p$	36	$\text{con } t$	108	Θ	146
$\text{ante } P$	36	\mathcal{L}	108	Ξ	146
$\tau(Q; \mathcal{X}), \tau(Q), \tau_Q$	39	\mathcal{L}_{fin}	109	$\varphi^{\text{rd}}, X^{\text{rd}}$	146
$\mathbf{P}(Q), \mathbf{P}_Q$	48	$\text{Sf } \varphi$	110	ZFC_{fin}	155
$Q \vdash (X, \varphi, p)$	48	$\text{len } \varphi$	110	ZFC_{-}	155
$X \leftrightarrow X_0$	49	$\text{rk } \varphi$	110	ZFC	155

ZFC ₊	155	φ_{lim}	168	$E[t \mid X, s]$	179
L_-	162	$\mu_{t X}$	171	$P(t \in \underline{V} \mid X, s)$	179
L_{ZFC}	165	$E[t \mid X]$	171		
ω_b, ω_t	165	$\mu_{\mathbf{t} X}$	172		