Mathematical Logic II

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Preface

This compendium will be the curriculum text for the course “Mathematical Logic II” at the University of Oslo, Spring 2006. This compendium is a revised version of the compendium used for the same course in the Spring 2004. A few details are added, some typos are corrected, some exercises that were unsolvable has been rewritten and other minor details have been changed.

In its present form, this compendium may be used free of charge by anyone, but if someone uses it for an organized course, the author would like to be informed.

Blindern, 13-12-05

Dag Normann
Introduction

This compendium is written primarily as a text for the course Mathematical Logic II given at the University of Oslo, Norway. However, much more stuff is included than what is needed for this; we have included material that is of interest to the logic group in Oslo, and that the reader might like to learn by self study. The course Mathematical Logic II is a continuation of the course Mathematical Logic, and we assume that the reader has passed this course or is familiar with logic up to the same level. In the Appendix we have collected some propositions without proofs. These are precise versions of theorems assumed to be known to the reader.

We will assume that the reader is familiar with first order logic at an elementary level, including the soundness theorem, the deduction theorem and the completeness theorem for countable theories. We will assume knowledge of the Henkin-style proof of the completeness theorem. The theorem of constants will be used extensively. At some occasions we will use the full completeness theorem, though a full proof use Zorn’s Lemma or comparative means. Since most intended readers also will follow a course in axiomatic set theory, the verification of Zorn’s Lemma from the Axiom of Choice is left for that course.

This compendium is essentially consisting of two parts, Model Theory and Computability Theory. In model theory we mainly restrict ourselves to the classical topics like quantifier elimination and omitting/realizing types, but have included a brief introduction to finite model theory. In the computability theory part we use a Kleene-style introduction, and give an introduction to the recursion theorem, c.e. sets, Turing degrees, a basic priority argument and the existence of minimal degrees. We also include a chapter on generalized computability theory.

Now we will settle some terminology and conventions that we will use in the text. Our terminology and our conventions are chosen to be compatible with the beginners course in logic at the University of Oslo, but our terminology is sufficiently standard for this to cause no problems for readers with a different background.

- $L$ will be a first order language. Unless specified otherwise, $L$ will be a language with equality.
- $c_1, c_2$ etc. will be names (or constant symbols) in a formal language.
- $f, g, etc.$ will be function symbols in a formal language. Sometimes it will be convenient to consider a name as a function symbol of arity 0.
- $R, P, Q$ etc. will be relation symbols in a formal language.
- $v_1, v_2, \ldots$ is the formal list of variables, but we will mainly use $x, y$ etc. instead.
- $t_1, t_2$ etc. will denote terms.
- $\phi, \psi, \xi$, etc. will denote formulas or sentences.
A structure $\mathfrak{A}$ for a language $L$ will consist of a domain $A$ and interpretations $l^A$ for each name, function symbol and relation symbol $l$ ( $l$ for ‘letter’) in $L$. Whenever we use the Gothic letters $\mathfrak{A}$ and $\mathfrak{B}$, we will use the corresponding roman letters $A$ and $B$ for the domains of these structures without always mentioning it.

$s$ will denote an assignment function, a function assigning an element of a structure $\mathfrak{A}$ to each variable $v$. $\phi[s]$ will be used intuitively as the formula $\phi$ where all free variables are replaced by names for elements in $A$ according to the assignment $s$.

Technically we define

\[ c^\mathfrak{A}[s] = c^A, \]
\[ v^\mathfrak{A}_i[s] = s(v_i), \]
\[ (ft_1\ldots t_n)^\mathfrak{A}[s] = f^A(t_1^\mathfrak{A}[s],\ldots,t_n^\mathfrak{A}[s]). \]

A formula $\phi$ will technically be interpreted as a function $\phi^\mathfrak{A}$ from the set of assignments to the set $\mathbb{B} = \{\top, \bot\}$ of truth values. How this is done is well known from any beginners course in logic.

We will use the terminology $t_{x_1,\ldots,x_n}$ and $\phi_{t_1,\ldots,t_n}$ for simultaneous instantiations of the terms $t_1,\ldots,t_n$ for the free occurrences of the variables $x_1,\ldots,x_n$. In the latter case we assume without saying that the terms are substitutable for the variables in question.

Since the notation above is hard to read, we will sometimes write

\[ \phi(x_1,\ldots,x_n) \]

when $\phi$ is a formula with at most $x_1,\ldots,x_n$ free. When we use this notation, we may write

\[ \phi(t_1,\ldots,t_n) \]

instead of

\[ \phi_{t_1,\ldots,t_n}. \]

We have chosen to use the notation

\[ \{\cdot;\cdot\cdot\} \]

instead of the commonly used

\[ \{\cdot|\cdot\cdot\}, \]

because in some expressions we would otherwise have a multiple purpose use of $|$ that might cause confused readers.

Each chapter 1-4 is supplied with a set of exercises at the end, some simple and some hard. The exercises are integrated parts of the text, and at the end the students are assumed to have worked through most of them.

Chapter 5 will consist of just exercises. These are not integrated parts, but challenges the eager reader might like to face.
Chapter 1
Classical Model Theory

In model theory we investigate the connection between particular first order theories or first order theories of a special kind and the class of models for that theory. The completeness theorem is one of the most basic tools in model theory, and some of the applications will refer to the actual proof of the completeness theorem. We will also, without hesitation, use the axiom of choice.

1.1 Embeddings and isomorphisms

In this section we will discuss more systematically ways of comparing structures of a 1. order language $L$. Some of these concepts may be known to the reader, but for the sake of completeness we include all relevant definitions.

Definition 1.1.1 Let $L$ be a first order language, $\mathfrak{A}$ an $L$-structure with domain $A$.

A substructure $\mathfrak{B}$ with domain $B$ will be an $L$-structure such that

1. $B \subseteq A$.
2. $c^B = c^\mathfrak{A}$ whenever $c$ is a name (constant symbol) in $L$.
3. $f^B(a_1, \ldots, a_n) = f^\mathfrak{A}(a_1, \ldots, a_n)$ whenever $f$ is a function symbol in $L$ of arity $n$ and $a_1, \ldots, a_n \in B$.
4. $R^B = R^\mathfrak{A} \cap B^n$ whenever $R$ is a relation symbol in $L$ of arity $n$.

Definition 1.1.2 A first order theory $T$ is open if all non-logical axioms are open, i.e. have no quantifiers.

Lemma 1.1.3 Let $\mathfrak{A}$ be an $L$-structure.

a) Let $B \subseteq A$. Then $B$ is the domain of a substructure $\mathfrak{B}$ of $\mathfrak{A}$ if and only if $B$ is closed under all interpretations $f^\mathfrak{A}$ where $f$ is a function symbol (or a constant) in $L$. 
b) Let $\mathfrak{B}$ be a substructure of $\mathfrak{A}$. Let $\phi$ be an open $L$-formula and $s$ an assignment over $\mathfrak{B}$ (and then also over $\mathfrak{A}$.) Then $\phi[s]$ is true in $\mathfrak{B}$ if and only if it is true in $\mathfrak{A}$.

Proof
The proof of a) is trivial and the proof of b) is trivial by induction on the subformula relation.

Given a structure $\mathfrak{A}$, it is sufficient to consider the class of substructures for most purposes where this is relevant. However, sometimes we want to go the other way, we might like to extend a structure. This is done e.g. in moving from the natural numbers to the integers, further on to the rationals and maybe continuing to the algebraic numbers, the real numbers or the complex numbers. Whenever convenient we may consider a natural number as a complex number, though using a standard representation of complex numbers in set theory, there is a long way to go. In the extreme, a natural number will be a finite set totally ordered by the $\in$-relation, an integer will be an equivalence class of pairs of natural numbers, a rational number will be an equivalence class of pairs of integers, a real will be a Dedekind cut of rational numbers (or even worse, an equivalence class of Cauchy sequences of rational numbers) and a complex number will be an ordered pair of real numbers. Of course there is no mathematical gain in thinking of complex numbers as such monsters, it is easier to assume that $\mathbb{N} \subset \mathbb{Z} \subset \mathbb{R} \subset \mathbb{C}$ and just base our mathematics on our intuitive understanding of these sets.

When we in the context of a logic course want to extend a structure, it is standard procedure to use the completeness theorem in such a way that we ensure that a copy of the original structure will be a substructure. We have to do some set-theoretical hocus pocus in order to obtain an actual extension. We will now develop the concepts needed in order to make this precise:

**Definition 1.1.4** let $L$ be a first order language and let $\mathfrak{A}$ and $\mathfrak{B}$ be $L$-structures.

a) An embedding from $\mathfrak{B}$ to $\mathfrak{A}$ is a map $\pi : B \rightarrow A$ that is one to one and such that

- $\pi(c^\mathfrak{B}) = c^\mathfrak{A}$ for each name $c$.
- $\pi(f^\mathfrak{B}(a_1, \ldots, a_n)) = f^\mathfrak{A}(\pi(a_1), \ldots, \pi(a_n))$ for all $a_1, \ldots, a_n \in B$.
- $R^\mathfrak{B}(a_1, \ldots, a_n) \Leftrightarrow R^\mathfrak{A}(\pi(a_1), \ldots, \pi(a_n))$ for all $a_1, \ldots, a_n \in B$.

b) An isomorphism will be an embedding that is onto.

**Examples 1.1.5**

a) If we consider $\langle \mathbb{R}, 0, +, < \rangle$ and $\langle \mathbb{R}^+, 1, \cdot, < \rangle$ as two structures for the language with one name, one binary function and one binary relation, then $\pi(x) = e^x$ is an isomorphism.

b) If we consider reals as sets of Dedekind cuts of rationals with the inherited algebraic structure, then $\pi(q) = \{ r \in \mathbb{Q} : r < q \}$ will be an embedding of $\mathbb{Q}$ into $\mathbb{R}$. 

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c) If $\mathfrak{A} = (\mathbb{Z}, +, \cdot, 0, 1)$ is seen as a structure for the language of ring theory, then the set of non-negative integers with the inherited algebra will be a substructure (but not a subring) of $\mathfrak{A}$.

There are two traditions for defining the interpretation of a formula over a structure. One is using assignments as we did in the introduction. Then we define the truth value of any formula relative to a given assignment. The other tradition is to extend the language with names for each element in the structure, and then only interpret the sentences in this extended language. These approaches are equivalent, but it turns out that the latter approach is more easy to use for many of the purposes in model theory.

**Definition 1.1.6** Let $L$ be a first order language, $\mathfrak{A}$ an $L$-structure.

a) By $L(\mathfrak{A})$ we mean the language $L$ where we add one new constant symbol $c_a$ for each $a \in A$.

b) As a default, we will consider $\mathfrak{A}$ as an $L(\mathfrak{A})$-structure via the interpretation $c^\mathfrak{A}_a = a$.

See Exercise 1.3 for the tedious, but simple observations that have to be made. Recall that a literal is a formula that is either an atomic formula or the negation of an atomic formula.

**Definition 1.1.7** Let $L$ be a first order language and let $\mathfrak{A}$ be an $L$-structure. By the diagram of $\mathfrak{A}$, $D(\mathfrak{A})$, we mean the set of variable-free literals in $L(\mathfrak{A})$ that are true in $\mathfrak{A}$.

**Remark 1.1.8** In some textbooks, the diagram is defined as the theory where all variable free sentences in $L(\mathfrak{A})$ that are true in $\mathfrak{A}$ are axioms. This will be an equivalent theory to ours, so the choice of definition is a matter of taste.

The diagram will play an important part when we want to show that we have extensions of a structure $\mathfrak{A}$ with some desired properties; we just have to show that the diagram of $\mathfrak{A}$ is consistent with the desired properties:

**Lemma 1.1.9** Let $L$ be a first order language, $\mathfrak{A}$ and $\mathfrak{B}$ $L$-structures. Then the following are equivalent:

i) There is an embedding $\pi$ of $\mathfrak{A}$ into $\mathfrak{B}$.

ii) We may give an interpretation of the extra constants $c_a$ from $L(\mathfrak{A})$ in $\mathfrak{B}$ such that $\mathfrak{B}$ becomes a model for $D(\mathfrak{A})$.

**Proof**

Assume i). Let $(c_a)^{\mathfrak{B}} = \pi(a)$. With this interpretation, ii) is verified.

Assume ii). Let $\pi(a) = (c_a)^{\mathfrak{B}}$. Since $L$ is a language with equality, and $\mathfrak{B}$ is a model for the diagram of $\mathfrak{A}$, $\pi$ will be an embedding.
It will be a matter of taste if we talk about extensions of structures or of embeddings into other structures, see Exercise 1.4. Our first result is a typical model-theoretical one, we connect the logical form of the axioms in a theory with the relations between models and their substructures. Recall that two first order theories are equivalent if they have the same theorems. (A consequence will be that they are based on the same language; why?)

**Theorem 1.1.10** (Las-Tarski)

Let $T$ be a first order theory over a language $L$. Then the following are equivalent:

1. For all $L$-structures $\mathfrak{A}$ and $\mathfrak{B}$, if $\mathfrak{B}$ is a model for $T$ and $\mathfrak{A}$ is a substructure of $\mathfrak{B}$, then $\mathfrak{A}$ is a model for $T$.

2. $T$ is equivalent to an open theory $T'$.

**Proof**

By the completeness theorem, two theories will be equivalent if and only if they have the same models. Thus 2. $\implies$ 1. is already proved.

Now assume that 1. holds. Let the non-logical axioms in $T'$ be the open theorems of $T$. Then $T'$ is an open theory, and all models for $T$ will be models for $T'$. In order to prove that $T$ and $T'$ are equivalent, we must prove the converse. So, let $\mathfrak{A}$ be a model for $T'$. In order to prove that $T'$ is a model for $T$, consider the theory

$$T^* = \mathcal{D}(\mathfrak{A}) \cup T.$$

**Claim**

$T^*$ is consistent.

**Proof**

Assume not. Then there are $\phi_1, \ldots, \phi_n$ in the diagram of $\mathfrak{A}$ such that $T \cup \{\phi_1, \ldots, \phi_n\}$ is inconsistent, or alternatively, that $T \vdash \neg(\phi_1 \land \cdots \land \phi_n)$.

Each $\phi_i$ will be an instance of some literal $\psi_i$ in $L$, and we may let the same name $c_a$ be substituted for the same variable $z$ in all the literals under consideration. Thus there are variables $z_1, \ldots, z_m$, new names $c_{a_1}, \ldots, c_{a_m}$ and literals $\psi_1, \ldots, \psi_n$ in $L$ such that

$$\phi_i = (\psi_i)_{z_1,\ldots,z_m}^{c_{a_1},\ldots,c_{a_m}}$$

for each $i \leq n$.

The names $c_{a_i}$ will not occur in the non-logical axioms in $T$, so by the theorem of constants, $T \vdash \neg(\psi_1 \land \cdots \land \psi_n)$.

However, each $\phi_i$ is a literal that is true in $\mathfrak{A}$, so $\phi_1 \land \cdots \land \phi_n$ is a valid instance of $\psi_1 \land \cdots \land \psi_n$. This contradicts that $\neg(\psi_1 \land \cdots \land \psi_n)$ is valid in $\mathfrak{A}$, and the claim is proved.

Let $\mathfrak{B}$ be a model for $T^*$. There will be an embedding of $\mathfrak{A}$ into $\mathfrak{B}$. By Exercise 1.4, $\mathfrak{A}$ is isomorphic to a substructure $\mathfrak{A}'$ of $\mathfrak{B}$, which by the assumption is a model for $T$. Thus $\mathfrak{A}$ is a model for $T$.

This ends the proof of the theorem.
Sometimes in model theory we will construct a structure by piecewise extensions, quite often using the completeness theorem at each step. In order to handle such constructions we need the concept of directed limit, a concept you find in other parts of mathematics as well.

**Definition 1.1.11** Let \( \langle I, \prec \rangle \) be a partial ordering. The ordering is called directed if

\[
\forall i \in I \forall j \in I \exists k \in I (i < k \land j < k).
\]

A directed ordering is a generalization of an increasing sequence.

**Definition 1.1.12** Let \( L \) be a first order language.

a) A directed system of \( L \)-structures will consist of

1. A directed ordering \( \langle I, \prec \rangle \)
2. An \( L \)-structure \( A_i \) for each \( i \in I \)
3. An embedding \( \pi_{ij} : A_i \rightarrow A_j \) whenever \( i \leq j \) in \( I \)

such that \( \pi_{ii} \) is the identity function on each \( A_i \) and such that whenever \( i \leq j \leq k \), then \( \pi_{ik} = \pi_{jk} \circ \pi_{ij} \) (where \( \circ \) denotes composition).

b) Using the notation from a), a directed limit of the system will be an \( L \)-structure \( A \) and embeddings \( \pi_i : A_i \rightarrow A \) for each \( i \in I \) such that

- If \( i \leq j \), then \( \pi_i = \pi_j \circ \pi_{ij} \).
- Let \( B \) be an \( L \)-structure and for each \( i \in I \), let \( \eta_i \) be an embedding from \( A_i \) to \( B \) such that \( \eta_i = \eta_j \circ \pi_{ij} \) whenever \( i \leq j \).

Then there is a unique embedding \( \eta : A \rightarrow B \) such that \( \eta_i = \eta \circ \pi_i \) for each \( i \in I \).

**Remark 1.1.13** The concept of a directed limit is more general. In a category theoretical context it is often known as a colimit.

**Theorem 1.1.14** Each directed system

\[
\langle I, \prec, \{ A_i \}_{i \in I}, \{ \pi_{i,j} \}_{i \leq j} \rangle
\]

will have a directed limit, and this is unique up to isomorphisms.

**Proof**

The reader should verify for her/himself that if \( A \) and \( B \) both satisfy the requirements of a directed limit, then \( A \) and \( B \) are isomorphic, see Exercise 1.5.

Let \( X \) be the set of ordered pairs \((i, a)\) such that \( i \in I \) and \( a \in A_i \) where \( A_i \) is the domain of \( A_i \).

We let \((i, a) \approx (j, b)\) if for some \( k \in I \) we have that \( i \leq k, j \leq k \) and \( \pi_{ik}(a) = \pi_{jk}(b) \).
This definition does not depend on the choice of \( k \), and \( \approx \) will be an equivalence relation.

Let \( A \) be the set of equivalence classes \([a, i]\), and let \( \pi_i(a) = [(a, i)] \).

Whenever \( C_1, \ldots, C_n \) are equivalence classes, we may find an \( i \in I \) and \( a_1, \ldots, a_n \in A_i \) such that \( C_j = [(i, a_j)] \) for each \( j \leq n \). Here we use that \( (I, \prec) \) is a directed ordering.

We let

\[
f^A(C_1, \ldots, C_n) = [(i, f^{A_i}(a_1, \ldots, a_n))]
\]

and

\[
R^A(C_1, \ldots, C_n) \iff R^{A_i}(a_1, \ldots, a_n)
\]

whenever \( f \) and \( R \) are \( n \)-ary.

This definition does not depend on \( i \).

If \( B \) and \( \eta_i \) are as in Definition 1.1.12 b), then \( \eta \) defined by

\[
\eta([(a, i)]) = \eta_i(a)
\]

will be the unique embedding satisfying \( \eta_i = \eta \circ \pi_i \) for each \( i \in I \).

The details left out and the rest of the proof is left for the reader as Exercise 1.5.

### 1.2 Elementary embeddings

If two structures \( \mathfrak{A} \) and \( \mathfrak{B} \) for the same language \( L \) are isomorphic, they will have the same logical properties. In fact, this must be true as long as a logical property reflects the structure itself and not the underlying set-theoretical representation of the domains or of other parts of the interpretations. In particular, they will share the same properties expressible in higher order logic. Sometimes a concept is defined to be logical if it is invariant under isomorphisms. In the same spirit, we may call a concept algebraic if it is invariant under embeddability. It is not certain that all researchers in algebra will agree.

Concepts like finite, well ordering, countable and complete are logical in this sense, but they are not expressible in first order logic. In order to make a discussion of this precise, we need a new concept:

**Definition 1.2.1** Let \( L \) be a first order language, and let \( \mathfrak{A} \) and \( \mathfrak{B} \) be two \( L \)-structures.

We say that \( \mathfrak{A} \) and \( \mathfrak{B} \) are elementary equivalent if

\[
\mathfrak{A} \models \phi \iff \mathfrak{B} \models \phi
\]

whenever \( \phi \) is a sentence in \( L \).

Since it is useful to consider a structure as the same structure even if we interpret several additional names, we refer to \( L \) in the notation for elementary equivalence. We write

\[
\mathfrak{A} \equiv_L \mathfrak{B}
\]

when \( \mathfrak{A} \) and \( \mathfrak{B} \) are elementary equivalent with respect to sentences in \( L \).
In the beginning of the Twentieth century some mathematicians hoped that mathematics could be build on an even more solid ground by reducing it to logic. The Norwegian logician Thoralf Skolem demonstrated the limitations of this program.

**Theorem 1.2.2 (Skolem)**

Let $L$ be the first order language of number theory with constants for ‘zero’ and ‘one’, function symbols for ‘addition’, ‘multiplication’ and ‘exponentiation’ and a relation symbol for the ordering. Let

$$\mathfrak{N} = \langle \mathbb{N}, 0, 1, +, \cdot, \exp, < \rangle$$

be the standard model. Then there is another $L$-structure $\mathfrak{N}'$ that is elementary equivalent to $\mathfrak{N}$, but not isomorphic to $\mathfrak{N}$.

**Proof**

This is a consequence of the compactness theorem. Let $c$ be a new name, let $T$ be the set of sentences in $L$ valid in $\mathfrak{N}$, and let $T^*$ be $T$ extended by the axioms $k_n < c$ for each $n$, where $k_n$ is the numeral in $L$ denoting the number $n$.

Each finite subtheory of $T^*$ may use $\mathfrak{N}$ as a model interpreting $c$ as a natural number (depending on $T^*$). Then by the compactness theorem, there is a model $\mathfrak{N}'$ for $T^*$. This model will be elementary equivalent to $\mathfrak{N}$ because both are models for $T$, but they are clearly not isomorphic.

**Remark 1.2.3** The completeness theorem and the compactness theorem were not available to Skolem. His argument used another form of model theoretical construction, something we might call a reduced ultra-product construction, see Exercise 5.3.

By a similar argument we may show that the mathematically important concept of a well ordering is not definable in first order logic.

**Definition 1.2.4** Let $\langle X, < \rangle$ be a total ordering. $\langle X, < \rangle$ is a **well ordering** if every nonempty subset $Y \subseteq X$ has a least element.

**Theorem 1.2.5** There is no first order theory $T$ with the class of well orderings as its models.

**Proof**

The language of any such theory $T$ must be the one with just one binary relation symbol, we may as well denote it by $<$. A sentence $\phi$ in $L$ will be a theorem in $T$ if and only if $\phi$ is true for all well orderings. So let $T$ be the theory with all sentences $\phi$ true for all well orderings as its non-logical axioms. We will show that $T$ has a model that is not a well ordering.

Extend $L$ with new names $c_k$ for each natural number $k$ and let $T^*$ be $T$ extended with all axioms $c_{k+1} < c_k$ for $k \in \mathbb{N}$. 

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Each finite subtheory $T_0$ of $T^*$ will have a model, since there are arbitrarily long finite well orderings. By the compactness theorem, $T^*$ has a model

$$\mathfrak{A} = \langle A, <^A, \{c^A_k \mid k \in \mathbb{N}\} \rangle.$$ 

Clearly the set $\{c^A_k \mid k \in \mathbb{N}\}$ has no least element, so $\langle A, <^A \rangle$ is not a well ordering.

However

$$\langle A, <^\mathfrak{A} \rangle \models T,$$

so this structure has all the first order properties shared by all well orderings.

A concept of importance in logic in general and in model theory in particular is completeness.

**Definition 1.2.6** A first order theory $T$ over a language $L$ is complete if $T$ is consistent and $T \vdash \phi$ or $T \vdash \neg \phi$ for each sentence $\phi$ in $L$.

The following is a direct consequence of the completeness theorem:

**Lemma 1.2.7** Let $T$ be a first order theory over a language $L$. The following are equivalent:

1. $T$ is complete.
2. $\mathfrak{A} \equiv_L \mathfrak{B}$ for all models $\mathfrak{A}$ and $\mathfrak{B}$ for $T$.

If $L$ is a first order language, $\mathfrak{A}$ and $\mathfrak{B}$ are $L$-structures, and $\pi$ is an embedding from $\mathfrak{A}$ to $\mathfrak{B}$, we may view $\mathfrak{B}$ as a structure for $L(\mathfrak{A})$ using the interpretation

$$(c^A_a)^\mathfrak{B} = \pi(a).$$

An embedding of $\mathfrak{B}$ into $\mathfrak{A}$ is a 1-1 map $\pi$ that preserves the basic algebra of $\mathfrak{B}$. We will also be interested in maps that preserve the first order logic of $\mathfrak{A}$:

**Definition 1.2.8** Let $L$ be a first order language, and let $\mathfrak{A}$ and $\mathfrak{B}$ be two $L$-structures.

a) Assume that $\mathfrak{A}$ is a substructure of $\mathfrak{B}$. We say that $\mathfrak{A}$ is an elementary substructure of $\mathfrak{B}$ if $\mathfrak{A} \equiv_{L(\mathfrak{A})} \mathfrak{B}$.

b) Assume that $\pi$ is an embedding from $\mathfrak{A}$ to $\mathfrak{B}$ and view $\mathfrak{B}$ as a structure for $L(\mathfrak{A})$ as above.

We call $\pi$ an elementary embedding if $\mathfrak{A} \equiv_{L(\mathfrak{A})} \mathfrak{B}$.

We used the diagram $\mathcal{D}(\mathfrak{A})$ to axiomatize that $\mathfrak{A}$ essentially is a substructure. Likewise, we will need a formal theory that axiomatizes that $\mathfrak{A}$ essentially is an elementary substructure.
Definition 1.2.9 Let $L$ be a first order language, $\mathfrak{A}$ an $L$-structure. The theory of $\mathfrak{A}$, $Th(\mathfrak{A})$, will have the sentences in $L(\mathfrak{A})$ that are true in $\mathfrak{A}$ as the non-logical axioms.

We make the following observation:

Lemma 1.2.10 Let $L$ be a first order language, $\mathfrak{A}$ and $\mathfrak{B}$ two $L$-structures. The following are equivalent:

1. There is an elementary embedding $\pi$ from $\mathfrak{A}$ to $\mathfrak{B}$.

2. There is an interpretation of each new constant in $L(\mathfrak{A})$ as an element in $\mathfrak{B}$ such that $\mathfrak{B}$ becomes a model for $Th(\mathfrak{A})$.

Proof
Both directions are based on establishing the equality

$$\pi(a) = (c_a)^B.$$ 

1. $\Rightarrow$ 2.
Assuming 1., we interpret $c_a$ in $\mathfrak{B}$ by $(c_a)^B = \pi(a)$. Since $\pi$ is elementary, we have that $\mathfrak{A} \equiv L(\mathfrak{A}) \mathfrak{B}$, and thus $\mathfrak{B}$ will be a model for $Th(\mathfrak{A})$.

2. $\Rightarrow$ 1.
Assuming 2., we define $\pi$ by $\pi(a) = (c_a)^B$.
Let $\phi$ be a sentence in $L(\mathfrak{A})$. Since $\mathfrak{B} \models Th(\mathfrak{A})$ we must have that

$$\mathfrak{A} \models \phi \Rightarrow \mathfrak{B} \models \phi.$$

$$\mathfrak{A} \models \neg \phi \Rightarrow \mathfrak{B} \models \neg \phi.$$

It follows that $\pi$ is an embedding and that $\pi$ is elementary.

Theorem 1.2.11 Let $L$ be a first order theory, $\mathfrak{A}$ and $\mathfrak{B}$ two $L$-structures. Then the following are equivalent:

1. $\mathfrak{A} \equiv L \mathfrak{B}$.

2. There is an $L$-structure $\mathfrak{C}$ with elementary embeddings $\pi : \mathfrak{A} \rightarrow \mathfrak{C}$ and $\eta : \mathfrak{B} \rightarrow \mathfrak{C}$.

Proof
2. $\Rightarrow$ 1. is trivial, so assume that $\mathfrak{A}$ and $\mathfrak{B}$ are elementary equivalent. We will extend $L$ by adding names $c_a$ for each $a \in A$ and a disjoint set of names $d_b$ for $b \in B$.
Let $T = Th(\mathfrak{A}) \cup Th(\mathfrak{B})$. We will prove that $T$ is consistent.
Assume not. We will use that the set of axioms in $Th(\mathfrak{A})$ and the set of axioms in $Th(\mathfrak{B})$ are closed under $\wedge$.
Then there are sentences $\phi$ in $Th(\mathfrak{A})$ and $\psi$ in $Th(\mathfrak{B})$ such that $\vdash \neg(\phi \wedge \psi)$. 

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Let \( x_1, \ldots, x_n, y_1, \ldots, y_m \) be fresh and distinct variables and let \( \phi' \) and \( \psi' \) be \( L \)-formulas such that \( \phi \) is of the form
\[
(\phi')_{x_1, \ldots, x_n}
\]
and \( \psi \) is of the form
\[
(\psi')_{y_1, \ldots, y_m}.
\]
By the theorem of constants
\[
\vdash \neg (\phi' \land \psi'),
\]
and using propositional logic and rules for quantifiers
\[
\vdash \exists x_1 \cdots \exists x_n \phi' \rightarrow \forall y_1 \cdots \forall y_m \neg \psi',
\]
Now \( \exists x_1 \cdots \exists x_n \phi' \) will hold in \( \mathfrak{A} \) since \( a_1, \ldots, a_n \) makes \( \phi' \) true. Since \( \mathfrak{A} \) and \( \mathfrak{B} \) are elementary equivalent, \( \exists x_1 \cdots \exists x_n \phi' \) will be true in \( \mathfrak{B} \) as well.

As a consequence we must have that \( \forall y_1 \cdots \forall y_m \neg \psi' \) is true in \( \mathfrak{B} \). But this is impossible, since there is an instance that makes \( \psi' \) true in \( \mathfrak{B} \), namely \( \psi \). Thus we have observed a contradiction.

The assumption was that \( T \) is not consistent, and our conclusion will be that \( T \) is consistent and has a model \( \mathfrak{C} \). By Lemma 1.2.10 we see that 2. holds. This ends the proof of the theorem.

In Exercise 1.6 we give an application of this theorem.

In section 1.1 we discussed directed limits of directed systems. Such systems where all embeddings are elementary are of a particular interest:

**Definition 1.2.12** Let \( L \) be a first order language, \( \langle I, < \rangle \) be a directed set and let \( \langle \{ \mathfrak{A}_i \}_{i \in I}, \{ \pi_{ij} \}_{i < j} \rangle \) be a directed system of \( L \)-structures.

We call the system **elementary** if each \( \pi_{ij} \) is an elementary embedding.

The limit embeddings of an elementary system will be elementary, see Exercise 1.7.

We will prove two important theorems about elementary extensions, the so called Löwenheim-Skolem theorems. The theorems say that there is no way first order logic can distinguish between infinities that dominate the size of the language in question. One consequence will be the Skolem’s paradox;

**there is an elementary countable substructure of the reals.**

This is not a real paradox, it only demonstrates that first order logic is not adequate for proving, or even stating, that \( \mathbb{R} \) is uncountable.

**Theorem 1.2.13** Let \( \mathfrak{A} \) be an infinite first order structure over the language \( L \). Let \( X \) be any set.

Then there is an elementary extension \( \mathfrak{B} \) of \( \mathfrak{A} \) and an injective map \( \delta : X \to B \).

(We can find elementary extensions of \( \mathfrak{A} \) of arbitrarily large cardinality.)
Proof
Let $L^*$ be $L(A)$ extended with a name $d_b$ for each element $b \in X$. Let $T$ be the theory $Th(A)$ extended with the axioms $d_b \neq d_c$ whenever $b$ and $c$ are distinct elements of $X$.

Since $A$ is infinite, $A$ can be viewed as a model for each finite subtheory of $T^*$, so $T^*$ has a model $B$. Using Exercise 1.4 adjusted to elementary embeddings, we see that we may let $B$ be an extension of $A$. Let $\delta(b) = (d_b)^B$. $\delta$ will be injective, and the theorem is proved.

This was the easy upwards Löwenheim-Skolem theorem. We will now face the downwards Löwenheim-Skolem theorem. Since we will not bother to get involved in too much cardinality arithmetics, we restrict ourselves to the countable case, which is the most important and best known case.

**Theorem 1.2.14** Let $L$ be a countable first order language, and let $A$ be an $L$-structure. Then $A$ has a countable elementary substructure.

**Proof**
Let $\phi$ be a formula in $L$, $a_1, \ldots, a_n$ elements of $A$. In order to improve the readability, we will let

$$\phi(x_1, \ldots, x_n)$$

mean that all free variables in $\phi$ are among $x_1, \ldots, x_n$, and we will write

$$\phi(a_1, \ldots, a_n)$$

instead of

$$\phi_{x_1, \ldots, x_n}^{x_1, \ldots, x_n}.$$

Let $\phi(y, x_1, \ldots, x_n)$ be a formula in $L$. A Skolem function for $\phi$ is a function $h : A^n \to A$ such that whenever $a_1, \ldots, a_n$ are in $A$ then

$$\phi(h(a_1, \ldots, a_n), a_1, \ldots, a_n)$$

is true in $A$ whenever there is some $b \in A$ such that $\phi(b, a_1, \ldots, a_n)$ is true in $A$.

For each such formula $\phi$, we select one Skolem function $h_\phi$. If $n = 0$, $h_\phi$ will be a function of no variables, which will be just an element of $A$. In particular, $h_{y = y} \in A$. (We need the axiom of choice to justify the existence of Skolem functions and to select one for each formula $\phi$.)

If $f$ is a function symbol in $L$, $f^A$ will be one of the Skolem functions: Let $\phi$ be the formula $y = f(x_1, \ldots, x_n)$. Then $f^A = h_\phi$.

Since $L$ is countable, the set of Skolem functions will be countable.

Let $B_0$ be the empty set, and by recursion, let

$$B_{k+1} = \{ h(a_1, \ldots, a_n) ; h \text{ is one of the Skolem functions } h_\phi \text{ and } a_1, \ldots, a_n \in$$
By induction on $k$ we see that each $B_k$ will be countable, and that $B_k \subseteq B_{k+1}$. Moreover, $B_1 \neq \emptyset$ since $h_{y=y} \in B_1$.

Let $B$ be the union of the $B_k$’s. Then $B$ is a countable subset of $A$, and by construction, $B$ is closed under all the Skolem functions. By Lemma 1.1.3, $B$ is the domain of a substructure $\mathfrak{B}$ of $\mathfrak{A}$. It remains to prove that $\mathfrak{B}$ is an elementary substructure of $\mathfrak{A}$, which amounts to prove that whenever $a_1, \ldots, a_n$ are in $B$, then $\phi(a_1, \ldots, a_n)$ is true in $\mathfrak{B}$ if and only if it is true in $\mathfrak{A}$.

This is proved by induction on the complexity of $\phi$:
If $\phi$ is atomic, the equivalence follows from the fact that $\mathfrak{B}$ is a substructure of $\mathfrak{A}$.
If $\phi = \neg \psi$ or $\phi = \psi_1 \lor \psi_2$, the induction step is trivial.
If $\phi(a_1, \ldots, a_n) = \exists y \psi(y, a_1, \ldots, a_n)$, and $\phi(a_1, \ldots, a_n)$ is true in $\mathfrak{B}$ via some $b$, then, by the induction hypothesis, it will be true in $\mathfrak{A}$ via the same $b$.
If on the other hand $\phi(a_1, \ldots, a_n)$ is true in $\mathfrak{A}$ via some $b$, it is true in $\mathfrak{A}$ via $h_\psi(a_1, \ldots, a_n)$.
Since $h_\psi(a_1, \ldots, a_n) \in B$, we may use the induction hypothesis and see that $\phi(a_1, \ldots, a_n)$ is true in $\mathfrak{B}$ via $h_\psi(a_1, \ldots, a_n)$.

This ends the proof.

Readers familiar with cardinal arithmetics will observe that the same argument can be used to prove the general version:

**Theorem 1.2.15** Let $\kappa$ be an infinite cardinal number. Let $L$ be a first order language of cardinality $\leq \kappa$, and let $\mathfrak{A}$ be an $L$-structure. Then $\mathfrak{A}$ has an elementary substructure $\mathfrak{B}$ of cardinality at most $\kappa$.

### 1.3 Complete Theories

In the previous section, we defined a consistent first order theory to be complete if all sentences of the language can be proved or disproved. In this section we will look for criteria for a theory to be complete.

#### 1.3.1 Categorical theories

A first order theory is **categorical** if all models are isomorphic. If a theory is categorical, it is a consequence of the upwards Löwenheim-Skolem theorem that all models are finite. Thus this concept is of a rather limited interest. The following is more interesting:

**Definition 1.3.1** Let $L$ be a countable first order language, and $T$ a theory over $L$. $T$ is **$\omega$-categorical** if all countable models are isomorphic.

It is easy the see that a consistent $\omega$-categorical theory must be complete, see Exercise 1.8. In Exercise 1.12 we see that the converse is not true.
Example 1.3.2 Let $DO$ be the first order theory over the language $L$ with equality and one binary relation symbol $<$, and with the following axioms:

**DO-1** \( \neg(x < x) \)

**DO-2** \( x < y \land y < z \to x < z \)

**DO-3** \( x < y \lor y < x \lor x = y \)

**DO-4** \( x < y \to \exists u \exists v \exists w (u < x \land x < v \land v < y \land y < w) \)

The first three axioms tells us that $<$ is a total ordering, while DO-4 implies that there is no largest element, no least elements, and that the elements in the ordering are densely ordered. $DO$ stands for *dense ordering*.

We claim that $DO$ is $\omega$-categorical. We give an outline of the proof, and leave the details as Exercise 1.9

Let $\langle A, <_A \rangle$ and $\langle B, <_B \rangle$ be two countable models for $DO$. A *finite partial isomorphism* will be an order preserving map $p : K \to B$ where $K \subseteq A$ is finite. The set of finite partial isomorphisms have the following extension properties:

- If $p$ is a finite partial isomorphism and $a \in A$, then $p$ can be extended to a finite partial isomorphism $q$ defined on $a$.

- If $p$ is a finite partial isomorphism and $b \in B$, then $p$ can be extended to a finite partial isomorphism $q$ with $b$ in its range.

Given enumerations $A = \{a_n : n \in \mathbb{N}\}$ and $B = \{b_n : n \in \mathbb{N}\}$ we can construct an increasing sequence of finite partial isomorphisms $\{p_n\}_{n \in \mathbb{N}}$ securing that $p_{2n}(a_n)$ is defined and that $b_n$ is in the range of $p_{2n+1}$. The limit of these finite partial isomorphisms will be an isomorphism between the two structures.

We will discuss $\omega$-categoricity in more depth in the paragraph on element types and in the chapter on finite model theory. The method of proof is for obvious reasons called a *back-and-forth construction*. Sometimes in the literature, it is called a *zig-zag-construction*.

### 1.3.2 Elimination of Quantifiers

One of the key success stories of model theory is that of proving completeness of a theory via elimination of quantifiers, and thereby proving theorems about algebraic theories of genuine interest. Our key example will be field theory, but the main application is the theory of real closed fields. We have used traditional texts on model theory as sources for our exposition in this and later sections on model theory, in particular Sacks [2]. We will assume that the reader is familiar with field theory, but give a brief introduction for the sake of completeness.
Field Theory

The language of field theory that we will use, will consist of three constants, 0, 1 and \(-1\) and two binary function symbols \(+\) and \(\cdot\). As logicians, we should make a clear distinction between these symbols and the symbols we use for the particular interpretations, but unless forced by the circumstances, we will not do so.

Our axioms for field theory, \(F\), will be:

\[
\begin{align*}
\text{F 01} & \quad (x + y) + z = x + (y + z) \\
\text{F 02} & \quad x + 0 = x \\
\text{F 03} & \quad x + (-1 \cdot x) = 0 \\
\text{F 04} & \quad x + y = y + x \\
\text{F 05} & \quad (x \cdot y) \cdot z = x \cdot (y \cdot z) \\
\text{F 06} & \quad x \cdot 1 = x \\
\text{F 07} & \quad x \neq 0 \rightarrow \exists y(x \cdot y = 1) \\
\text{F 08} & \quad x \cdot y = y \cdot x \\
\text{F 09} & \quad x \cdot (y + z) = x \cdot y + x \cdot z \\
\text{F 10} & \quad 0 \neq 1
\end{align*}
\]

We will use the standard algebraic notation like e.g. \(xy\) and \(-x\) for \(x \cdot y\) and \((-1) \cdot y\) resp.

A field will be a model for this theory. Each natural number \(n\) will have an interpretation in a field \(\mathfrak{F}\). If \(\mathfrak{F}\) is a field and \(p\) is a prime, we say that the field have characteristic \(p\) if \(\mathfrak{F}||= p = 0\).

\(\mathbb{Z}_p = \{0, \ldots, p - 1\}\) where addition and multiplication are carried out modulo \(p\) (where \(p\) is a prime) is an example of a field of characteristic \(p\). If the reader is not familiar with this fact she/he should verify it by her/himself. A field will be of characteristic 0 if it is not of characteristic \(p\) for any prime \(p\) (a non-prime cannot be the characteristic of a field). The distinction between finite and infinite characteristics cannot be made in first order logic, neither can the distinction between finite and infinite fields. See Exercise 1.10 for more precise statements about this.

Other familiar examples of fields are \(\mathbb{Q}\), \(\mathbb{R}\) and \(\mathbb{C}\) with the standard algebraic structures. All these fields will have characteristic 0.

**Definition 1.3.3** A field \(\mathfrak{F}\) is *algebraically closed* if each polynomial \(P\) of one variable and degree \(> 0\) will have a root. This can be expressed by the infinite set of axioms

\[y_n \neq 0 \rightarrow \exists x(y_n x^n + y_{n-1} x^{n-1} + \cdots + y_0 = 0),\]
where $x^n$ has its usual mathematical meaning, and $n \geq 1$.
We let $ACF$ denote the theory of algebraically closed fields, and $ACF(p)$ or $ACF(0)$ denote the extension where the characteristic of the field is specified.

**Lemma 1.3.4** An algebraically closed field is never finite.

**Proof**
Let $a_1, \ldots, a_n$ be elements of an algebraically closed field, and let

$$P(x) = (x - a_1) \cdots (x - a_n) - 1.$$ 

If $P(x) = 0$, then $x \neq a_1, \ldots, a_n$.

We will prove that $ACF(0)$ and $ACF(p)$ for each $p$ are complete, combining methods from logic and basic algebra. In proving this, we will isolate concepts of a more general interest, and we will state and prove general lemmas and theorems that will have the completeness of these theories as consequences. At the end, we will introduce the theory of Real Closed Fields, and use our method to prove completeness of this theory as well.

**The isomorphism property**

**Definition 1.3.5** Let $T$ be a first order theory.
We say that $T$ has the isomorphism property if whenever $\mathcal{A}$ and $\mathcal{B}$ are models for $T$, $\mathcal{A}_0$ and $\mathcal{B}_0$ are substructures of $\mathcal{A}$ and $\mathcal{B}$ resp. and $\pi_0$ is an isomorphism from $\mathcal{A}_0$ to $\mathcal{B}_0$, then $\pi_0$ can be extended to an isomorphism between submodels $\mathcal{A}_1$ and $\mathcal{B}_1$ of $\mathcal{A}$ and $\mathcal{B}$.

**Example 1.3.6** The theory $DO$ has the isomorphism property:

Let $\langle A, < \rangle$ be an ordering. A gap in $A$ will be an empty open interval $(a, b)$, $(-, a)$ or $(b, -)$ induced by a pair $a < b$ with no $c$ between them, a minimal $a$ or a maximal $b$. The gap structure reflects all instances of failure of the DO-axioms.

If we have two dense orderings and two isomorphic suborderings, either these suborderings are themselves dense, or they will have isomorphic gap structures.

Using the proof of the $\omega$-categoricity of $DO$ we see that each pair of gaps can be filled by copies of $\mathbb{Q}$ taken from the given dense orderings. The result will be isomorphic submodels. We leave the details for the reader.

**Example 1.3.7** The theory of fields has the isomorphism property:

A substructure of a field will contain 0, 1 and $-1$, and will be closed under summation and multiplication, so a substructure is actually a subring. Thus the task is to show that isomorphic subrings can be extended to isomorphic subfields inside the given fields. This is basic algebra, we have isomorphic "quotient-structures" of the two rings, and interpreting each quotient in the fields will give isomorphic subfields.

We also use standard algebra to prove
Lemma 1.3.8 The theory $ACF$ will have the isomorphism property

Proof
Since this compendium is not a textbook in algebra, we will not give all the details.

Given a field $\mathbb{F}_0$ and an irreducible polynomial $P(x)$ over that field we may extend the field by a root of that polynomial in a purely algebraic way, by considering the field of all polynomials $Q(x)$ modulo $P(x)$. If $\mathbb{F}_0$ is a subfield of an algebraically closed field $\mathbb{F}$, and $a \in \mathbb{F}$ is a root of $P$, then the map $Q(x) \mapsto Q(a)$ is an isomorphism between the formal extension and a larger subfield of $\mathbb{F}$. (Euclid’s algorithm for the largest common divisor applied to polynomial division is central here).

The algebraic closure of a field can be described as the result of a transfinite sequence of such formal extensions, so the algebraic closures of isomorphic fields within algebraically closed fields will be isomorphic.

In Exercise 1.11 we will give an alternative proof.

Lemma 1.3.9 Let $T$ be a first order theory over a language $L$ and assume that $T$ has the isomorphism property.

Let $T'$ be obtained from $T$ by adding new constants to $L$, but no new non-logical axioms. Then $T'$ has the isomorphism property.

The proof is trivial and is left for the reader.

The submodel property

Definition 1.3.10 Let $\phi$ be a formula in a first order language $L$. We call $\phi$ simple if $\phi$ is of the form $\exists x \psi$ where $\psi$ is open.

Definition 1.3.11 Let $L$ be a first order language, $T$ a theory over the language $L$.

We say that $T$ has the submodel property if whenever $\mathcal{B}$ is a model of $T$, $\mathcal{A}$ is a submodel of $\mathcal{B}$ and $\phi$ is a simple sentence in $L(\mathcal{A})$, then $\phi$ is true in $\mathcal{A}$ if and only if it is true in $\mathcal{B}$.

Example 1.3.12 The theory $DO$ has the submodel property.

Dropping formalities, let $\phi$ be the sentence $\exists x \psi(x, a_1, \ldots, a_n)$. $\psi$ can only express that $x$ equals some $a_i$ or that it is larger than some $a_i$'s and smaller than others. If one such $x$ can be found in some dense ordering extending $\{a_1, \ldots, a_n\}$ then it can be found in all such dense orderings.

Example 1.3.13 Field theory does not have the submodel property

Consider the statement $\exists x(x \cdot x = 2)$. This sentence is not true for the field of rationals, but for the field of reals.

Lemma 1.3.14 The theory of algebraically closed fields has the submodel property.
Let $B$ be an algebraically closed field, and let $A$ be an algebraically closed subfield. Let $\phi = \exists x \psi$ be a simple sentence in $L(A)$. Then $\psi$ is a Boolean combination of atomic formulas $\xi(x, a_1, \ldots, a_n)$, where $a_1, \ldots, a_n$ are in $A$. The function symbols are denoting plus and times, so each term $t(x)$ in $L(A)$ will be equivalent to a polynomial in the variable $x$ with coefficients from $A$. Since $=$ is the only predicate, each atomic formula $\xi(x, a_1, \ldots, a_n)$ is equivalent to a polynomial equation $P_\xi(x) = 0$, where $P_\xi$ has coefficients from $A$. $P_\xi$ will have a root in $A$ if and only if it has a root in $B$.

If $\phi$ is true in $A$, then $\phi$ is trivially true in $B$. On the other hand, assume that $\phi$ is true in $B$ via $b$, i.e. $\psi(b, a_1, \ldots, a_n)$ is true in $B$. If one of the atomic subformulas $\xi$ are true for $x = b$, then $b$ is the root of the polynomial $P_\xi$ with coefficients from $A$, so $b$ is in $A$. On the other hand, if none of the subformulas $\xi$ of $\psi$ are true for $x = b$, we use that $A$ is infinite (see Lemma 1.3.4), and that a finite set of polynomials only will have a finite set of roots. Thus let $x = a$, where $a$ is in $A$ and $a$ is not making any of the subformulas $\xi$ of $\psi$ true. Then $a$ and $b$ satisfy the same atomic subformulas of $\psi$, so $\psi(a, a_1, \ldots, a_n)$ must be true. This means that $\phi$ is true in $A$.

Lemma 1.3.15 Let $L$ be a first order language and let $T$ be a theory over $L$ that has the submodel property. Let $T'$ be obtained from $T$ by adding new names to $L$ but no new non-logical axioms. Then $T'$ has the submodel property.

The proof is easy and is left for the reader.

Elimination of quantifiers

Definition 1.3.16 Let $L$ be a first order language, $T$ a theory over $L$. We say that $T$ accepts elimination of quantifiers if we for each formula $\phi$ have an open formula $\psi$ such that $T \vdash \phi \iff \psi$.

Lemma 1.3.17 Let $L$ be a first order language and $T$ a theory over $L$. Then $T$ accepts elimination of quantifiers if and only if each simple formula in $L$ is equivalent in $T$ to an open formula.

Proof
The ‘only if’-direction is obvious, so assume that each simple formula in $L$ will be equivalent in $T$ to an open formula.
By induction on the complexity of a formula $\phi$, we prove that it is equivalent in $T$ to an open formula, where we use the assumption to handle the quantifier case.
Lemma 1.3.18  Let $L$ be a first order language, $T$ a theory over $L$. Let $\phi$ be a sentence in $L$, and assume that whenever $\mathfrak{A}$ and $\mathfrak{B}$ are models for $T$ satisfying the same variable free sentences, then $\phi$ has the same truth values in $\mathfrak{A}$ and in $\mathfrak{B}$. Then $\phi$ is equivalent in $T$ to a variable free sentence in $L$.

Note
If a formula $\phi$ is variable free, it contains neither bound nor free variables, and thus no quantifiers. Another way of describing such formulas are as open sentences.

Proof
Let $T_0$ be $T$ extended with all variable free theorems in $T \cup \{ \phi \}$. It is sufficient to show that $\phi$ is a theorem in $T_0$.
Assume not, and let $\mathfrak{A}$ be a model for $T_0 \cup \{ \neg \phi \}$.
Let $\Delta$ be the set of variable free sentences true in $\mathfrak{A}$, and let $\mathfrak{B}$ be a model for $T \cup \Delta$. Then the same variable free sentences are true in $\mathfrak{A}$ and in $\mathfrak{B}$ (namely $\Delta$), so $\neg \phi$ will hold in $\mathfrak{B}$ as well.
$\mathfrak{B}$ was an arbitrary model for $T \cup \Delta$, so by the completeness theorem,

$$T \cup \Delta \vdash \neg \phi.$$ 

Then there are $\xi_1, \ldots, \xi_n \in \Delta$ such that

$$T, \phi \vdash \neg(\xi_1 \land \cdots \land \xi_n).$$

Then $\neg(\xi_1 \land \cdots \land \xi_n)$ is an axiom in $T_0$ and

$$\mathfrak{A} \models \neg(\xi_1 \land \cdots \land \xi_n).$$

This contradicts that $\mathfrak{A} \models \xi_i$ for each $i$.
The assumption was that $\phi$ is not a theorem in $T_0$. This led to a contradiction, so we are through.

Theorem 1.3.19 Let $L$ be a first order language, $T$ a theory over $L$ that has both the isomorphism property and the submodel property. Then $T$ accepts elimination of quantifiers.

Proof
By Lemma 1.3.17 it is sufficient to show that if $\phi = \exists x \psi(x, x_1, \ldots, x_n)$ where $\psi$ is open, then $\phi$ is equivalent in $T$ to an open formula. We may without loss of generality assume that $n \geq 1$. Extend $T$ to $T'$ by adding names $e_1, \ldots, e_n$, but no new non-logical axioms. By Lemmas 1.3.9 and 1.3.15, $T'$ has the isomorphism and submodel properties.
We will show that $\exists x \psi(x, e_1, \ldots, e_n)$ is equivalent in $T'$ to a variable free sentence $\xi(e_1, \ldots, e_n)$. Having done this, we conclude, using the theorem of constants, that $\exists x \psi(x, x_1, \ldots, x_n)$ is equivalent in $T$ to $\xi(x_1, \ldots, x_n)$.

By Lemma 1.3.18 it is sufficient to show that whenever $\mathfrak{A}$ and $\mathfrak{B}$ are two models
for $T'$ satisfying the same variable free formulas, then $\exists x \psi(x, e_1, \ldots, e_n)$ will have the same truth values in $\mathfrak{A}$ and $\mathfrak{B}$.

If $\mathfrak{A}$ and $\mathfrak{B}$ satisfy the same variable free formulas, the respective minimal substructures $\mathfrak{A}_0$ and $\mathfrak{B}_0$ consisting of all interpretations of closed terms will be isomorphic. By the isomorphism property for $T'$, this isomorphism can be extended to an isomorphism between two submodels $\mathfrak{A}_1$ and $\mathfrak{B}_1$ of $\mathfrak{A}$ and $\mathfrak{B}$ resp. Since $T'$ has the submodel property, we have that $\exists x \psi(x, e_1, \ldots, e_n)$ will have the same truth values in $\mathfrak{A}$ and $\mathfrak{A}_1$ and in $\mathfrak{B}$ and $\mathfrak{B}_1$. Since $\mathfrak{A}_1$ and $\mathfrak{B}_1$ are isomorphic, we are through.

Prime Models

**Definition 1.3.20** Let $L$ be a first order language with at least one constant symbol. Let $T$ be a first order theory over $L$, and let $\mathfrak{A}$ be an $L$-structure. $\mathfrak{A}$ is called a prime model for $T$ if $\mathfrak{A}$ can be embedded into any model $\mathfrak{B}$ for $T$.

Note: A prime model for $T$ is not necessarily a model for $T$. Actually we have

**Lemma 1.3.21** Let $L$ be a first order language, $T$ a theory over $L$ and $\mathfrak{A}$ a prime model for $T$. Then $\mathfrak{A}$ is a prime model for each extension $T'$ of $T$ over $L$.

The (trivial) proof is left for the reader.

**Example 1.3.22** Let $L$ be the language of $DO$ extended with the constants $c_i$ for $i \in \mathbb{N}$.

Let $DO^+$ be the theory $DO$ extended with the axioms $c_i < c_{i+1}$.

Then $\mathbb{N}$ with its usual ordering, and with $i$ as the interpretation of $c_i$, is a prime model for $DO^+$.

**Examples 1.3.23** The theories $ACF(0)$ and $ACF(p)$ have prime models $\mathbb{N}$ and $\mathbb{Z}_p$ respectively. These are also prime models for the theories of fields with a fixed characteristic.

In Exercise 1.12 We will see that $DO^+$ is a complete theory that is not $\omega$-categorical.

**Lemma 1.3.24** Let $T$ be a first order theory over a language $L$, and assume that $L$ has at least one constant symbol. Then the following are equivalent:

1. $T$ has a prime model.
2. Each variable free formula in $L$ is decidable in $T$.

**Proof**

First assume 1. Let $\mathfrak{A}$ be a prime model for $T$. Then for each model $\mathfrak{B}$ for $T$ and each variable free $\phi$ (the constant symbol in $L$ ensures that there are variable free formulas), the truth value of $\phi$ in $\mathfrak{B}$ will be the same as in $\mathfrak{A}$. Thus the existence of a prime model ensures that each variable free formula either is
true in all models for $T$, and thus is a theorem, or is false in all models, and thus is disprovable in $T$. Thus they are all decidable in $T$, and 2. is proved.

Now assume 2. For a model $\mathfrak{A}$ of $T$, let $\mathfrak{A}_0$ be the submodel consisting of all interpretations of closed terms $t$ in $\mathfrak{A}$.

If $\mathfrak{A}$ and $\mathfrak{B}$ are two models for $T$, we define $\pi : \mathfrak{A}_0 \rightarrow \mathfrak{B}_0$ by $\pi(t^\mathfrak{A}) = t^\mathfrak{B}$. Since all variable free formulas are decidable, $\pi$ is well defined and an isomorphism. This shows that $\mathfrak{A}_0$ is a prime model for $\mathfrak{A}$ whenever $\mathfrak{A} \models T$.

A completeness criterion

**Theorem 1.3.25** Let $L$ be a first order language and let $T$ be a theory over $L$ such that

1. $T$ is consistent.
2. The language $L$ of $T$ has at least one constant symbol.
3. $T$ has a prime model.
4. $T$ has the isomorphism property.
5. $T$ has the submodel property.

Then $T$ is complete.

**Proof**

By Theorem 1.3.19 we have that $T$ accepts elimination of quantifiers.

Let $\phi$ be a sentence in $L$ and let $\psi$ be open such that $T \vdash \phi \leftrightarrow \psi$. We may assume that $\phi$ is variable free, since otherwise we may replace it by $\phi(c, \ldots, c)$, where $c$ is a constant of the language.

Since $T$ has a prime model, $\psi$ is decidable in $T$ by Lemma 1.3.24, so $\phi$ is decidable in $T$, meaning that $T \vdash \phi$ or $T \vdash \neg \phi$.

**Corollary 1.3.26** The theories $ACF(0)$ and $ACF(p)$ for prime numbers $p$ are complete.

### 1.3.3 Real Closed Fields

One of the first success stories in model theory was the proof due to Tarski that the theory of Real Closed Fields permits elimination of quantifiers, and thus that it is complete.

**Definition 1.3.27** Let the theory $OF$ of ordered fields be field theory extended with one binary relation $<$, the axioms for total orderings and the axioms

\[
x < y \rightarrow x + z < y + z.
\]

\[
0 < x \land 0 < y \rightarrow 0 < xy.
\]
There are obvious models for $OF$, see Exercise 1.16. Working with ordered fields we may use most of our intuition about subfields of $\mathbb{R}$, though there will be ordered fields that are not subfields of $\mathbb{R}$. Let us consider an example:

**Example 1.3.28** Let $\mathfrak{Q}$ be the standard ordered field of rationals. Then each natural number $n$ is the interpretation of some term $\bar{n}$ in the language.

Let $c$ be a new constant, let $T = Th(\mathfrak{Q})$ and let $T^*$ be $T$ extended with the axioms $\bar{n} < c$ for each number $n$.

By the compactness theorem, $T^*$ has a model $\mathfrak{Q}^*$. This will be an ordered field, but certainly not a subfield of $\mathbb{R}$.

**Lemma 1.3.29** $OF \vdash 0 \neq x \rightarrow x^2 > 0$.

**Proof**

By the first new axiom it follows that $x < 0 \rightarrow -x > 0$. Since $(-x)^2 = x^2$ for any field, we may split the argument in the two cases $x > 0$ and $x < 0$ and use the second new axiom.

We will use standard algebraic properties of ordered fields without proving them in detail from the axioms. The uneasy reader should fill in the details for her/himself. We will show a little more care about facts that are proved using methods from analysis.

Recall from basic analysis that an ordering is complete if each bounded nonempty set has a least upper bound,

**Lemma 1.3.30** Each ordered field $\mathfrak{F}$ can be extended to an ordered field $\mathfrak{G}$ where the ordering is complete.

**Proof**

Let $\mathfrak{F}$ be an ordered field with domain $F$. Let $a < b$ be in $F$, then

$$a - 1 < a < \frac{a + b}{2} < b < b + 1,$$

so the ordering of $F$ will be a model for $DO$.

In a dense ordering $X$, a Dedekind cut will be a set $Y \subseteq X$ such that

$Y$ is bounded in $X$.

$Y$ is an initial segment of $X$.

$Y$ has no largest element.

The Dedekind cuts are ordered by inclusion. The original ordering may be embedded into the set of Dedekind cuts by

$$\pi(x) = \{y : y < x\}.$$

If $\mathfrak{F}$ is an ordered field, we may extend the algebra on $\mathfrak{F}$ to an algebra on the set of Dedekind cuts as follows
1. $X + X' = \{a - b : a \in X \land b \in X'\}$.

2. If $X$ and $Y$ both contain positive elements, let
   
   $$X \cdot Y = \{z \in F : \exists x \in X \exists y \in Y(x > 0 \land y > 0 \land z \leq x \cdot y)\}.$$

3. $(-1) \cdot X = \{-x : x \text{ is an upper bound for } X, \text{ but not a least upper bound for } X\}$.

4. For other pairs $X$ and $Y$ of Dedekind cuts we compute $X \cdot Y$ using 2. and 3.

The set of Dedekind cuts with the induced algebra is called the topological completion of $\mathbb{R}$. The topological completion will itself be a field, and $\mathbb{R}$ will be isomorphic to a subfield of the topological completion. The details are left as Exercise 1.17.

For an ordered field $\mathbb{F}$ we may define continuity of functions $f : F \to F$ by the standard $\epsilon - \delta$-definition. Then every polynomial $P(x)$ will define a continuous function. The details are left as Exercise 1.18.

**Lemma 1.3.31** Let $\mathbb{F}$ be an ordered field, $P(x)$ a polynomial over $\mathbb{F}$ of uneven degree $2n + 1$. Then $P(x)$ has a root in the topological completion.

**Proof**
The proof of this theorem for the reals, using the intermediate value theorem, works without change for any topologically closed ordered field, and thus for the topological closure of $\mathbb{R}$.

We also have

**Lemma 1.3.32** Let $a > 0$ in an ordered field. Then $a$ has a square root in the topological completion.

**Proof**
Let $C = \{b : b \leq 0 \lor (0 < b \land b^2 < a)\}$.

This is a Dedekind cut corresponding to $\sqrt{a}$.

**Definition 1.3.33** The theory $RCF$ of Real Closed Fields is obtained from $OF$ by adding the following axioms:

- $\exists y(x > 0 \rightarrow y^2 = x)$.
- $y_{2n+1} \neq 0 \rightarrow \exists x(y_0 + y_1 x + \cdots + y_{2n+1} x^{2n+1} = 0)$ for each number $n$.

Since the topological completion of the topological completion of an ordered field $\mathbb{F}$ is isomorphic to just one topological completion of $\mathbb{F}$, we just have shown that the topological completion of an ordered field will be a real closed field. The algebraic completion will be the least real closed subfield of the topological completion containing the image of the original $\mathbb{F}$ under the embedding. All the roots in the topological completion of a polynomial $P(x)$ in $\mathbb{F}$ will be in the algebraic completion. This is basic algebra.

The isomorphism property for $RCF$ is a consequence of the following...
Lemma 1.3.34 Let \( \mathfrak{F} \) and \( \mathfrak{G} \) be ordered fields where \( \mathfrak{F} \) is a subfield of \( \mathfrak{G} \) and \( \mathfrak{G} \) is real closed. Then the algebraic closure of \( \mathfrak{F} \) is isomorphic to a subfield of \( \mathfrak{G} \).

Proof
Let \( \mathfrak{F} \) and \( \mathfrak{G} \) be the topological completions of \( \mathfrak{F} \) and \( \mathfrak{G} \) resp. Since \( \mathfrak{F} \) is a subfield of \( \mathfrak{G} \), there is an inherited embedding \( \pi: \mathfrak{F} \rightarrow \mathfrak{G} \). \( \mathfrak{G} \) is isomorphic to a real closed subfield \( \mathfrak{G}' \) of \( \mathfrak{G} \) and \( \mathfrak{F} \) is isomorphic to a subfield \( \mathfrak{F}' \) of \( \mathfrak{F} \). Then \( \mathfrak{F}' \) is a subfield of \( \pi^{-1}(\mathfrak{G}') \), which is real closed, and by our construction of the algebraic completion, the algebraic completion of \( \mathfrak{F} \) is isomorphic to a subfield of \( \pi^{-1}(\mathfrak{G}') \) and thus to a subfield of \( \mathfrak{G} \).

Lemma 1.3.35 RCF has the submodel property.

Proof
Let \( \mathfrak{F} \) and \( \mathfrak{G} \) be two real closed fields, \( \mathfrak{F} \) a subfield of \( \mathfrak{G} \). Let \( \phi(x) \) be a quantifier free formula in \( L[\mathfrak{F}] \) and assume that \( \mathfrak{G} \models \exists x \phi(x) \). Each atomic formula \( \psi(x) \) in \( L[\mathfrak{F}] \) will be equivalent either to a polynomial equation \( P(x) = 0 \) or to a polynomial inequation \( P(x) > 0 \). Let \( \psi_1(x), \ldots, \psi_k(x) \) be the atomic subformulas of \( \phi \) and let \( P_1(x), \ldots, P_k(x) \) be the corresponding polynomials.

Let \( \{a_1, \ldots, a_n\} \) be the finite set of roots of the polynomials \( P_1(x), \ldots, P_k(x) \) in increasing order. The complements in \( \mathfrak{F} \) and in \( \mathfrak{G} \) will be corresponding finite unions of open intervals on which each of the polynomials \( P_i(x) \) will have a constant signature ‘positive’ or ‘negative’ (if not \( P_i(x) \) will have a root in the interval in the topological completion, and thus in the algebraic completion).

Let \( b \in |\mathfrak{G}| \) be such that \( \mathfrak{G} \models \phi(b) \). If \( b \) is a root in one of the polynomials, \( b \in |\mathfrak{F}| \) and we are through. If \( b \) is in one of the intervals, find \( a \in |\mathfrak{F}| \) in the corresponding interval. Then \( a \) and \( b \) will make exactly the same atomic subformulas of \( \phi \) true, so \( \phi(a) \Leftrightarrow \phi(b) \), and we are through.

We may now conclude

Theorem 1.3.36 The theory RCF is consistent, has a prime model and permits elimination of quantifiers. Consequently RCF is complete.

1.4 Element types

Throughout this section, we will let \( L \) be a fixed countable first order language, and we will let \( T \) be a complete theory over the language \( L \).

Definition 1.4.1 An \( n \)-type is a maximal set \( X \) of formulas \( \phi(x_1, \ldots, x_n) \) such that for all \( \phi_1, \ldots, \phi_k \) in \( X \) we have that

\[
T \vdash \exists x_1 \cdots \exists x_n (\phi_1 \land \cdots \land \phi_k).
\]
We will use the term complete for this property. Each consistent set of formulas with at most \( x_1, \ldots, x_n \) free can be extended to an \( n \)-type.

An \( n \)-type will be closed under \( \land \), and for each formula \( \phi \) with at most \( x_1, \ldots, x_n \) free the maximality of \( X \) will ensure that \( \phi \in X \) or \( \neg \phi \in X \).

**Definition 1.4.2** Let \( X \) be an \( n \)-type and \( A \) a model for \( T \). We say that \( A \) realizes \( X \) if for some \( a_1, \ldots, a_n \) in \( A \), \( A \models \phi(a_1, \ldots, a_n) \) for each \( \phi \in X \).

If \( A \) does not realize \( X \) we say that \( A \) omits \( X \).

**Example 1.4.3** Let \( \mathfrak{N} \) be the standard model for number theory and \( T = Th(\mathfrak{N}) \). Let \( X \) be a complete extension of \( \{k_n < x : n \in \mathbb{N}\} \) where \( k_n \) is the numeral for \( n \) in the language. This is a 1-type that is omitted by the standard model. The construction of a nonstandard elementary extension of \( \mathfrak{N} \) used to prove Theorem 1.2.2 amounts to constructing a model realizing this type.

The method used to prove Theorem 1.2.2 is quite general, and can be used to prove the following

**Theorem 1.4.4** Let \( X \) be an \( n \)-type. Then there is a model \( A \) for \( T \) realizing \( X \).

**Proof**
The proof is left as Exercise 1.20. Add new names for the \( n \) variables and use the compactness theorem.

Some \( n \)-types will be realized by every model of \( T \).

**Definition 1.4.5** An \( n \)-type \( X \) is principal if for some \( \phi \in X \) we have that

\[ T \vdash \phi \rightarrow \psi \]

for all \( \psi \in X \). We call \( \phi \) a generating formula. If \( X \) is not principal, we will use the term non-principal.

**Lemma 1.4.6** If \( X \) is a principal type, then \( X \) is realized in every model \( A \) of \( T \).

**Proof**
Any interpretation of \( x_1, \ldots, x_n \) that satisfies the generating formula will realize the type.

**Lemma 1.4.7** Let \( n \) be fixed.

Then the following are equivalent:

1. There are infinitely many \( n \)-types.
2. There is a non-principal \( n \)-type.
Proof
First assume that there is a non-principal \( n \)-type \( X = \{ \phi_i ; i \in \mathbb{N} \} \). For each \( i \in \mathbb{N} \) there must be a \( j_i > i \) such that
\[
T \not\models \phi_1 \land \cdots \land \phi_i \rightarrow \phi_{j_i}
\]
since otherwise the type will be principal. Thus, \( \{ \phi_1, \ldots, \phi_i, \neg \phi_{j_i} \} \) must be consistent, and there is an \( n \)-type \( X_i \) extending this set. Clearly, if \( j \geq j_i \), then \( X_i \) and \( X_j \) are different, so there are infinitely many \( n \)-types.

Now assume that all types are principal. Since \( L \) is countable, we will only be able to use countably many generating formulas, so the set of \( n \)-types is at most countable.

With the aim of obtaining a contradiction, assume that this set is infinite. Let \( \{ X_i \}_{i \in \mathbb{N}} \) be an enumeration of all the \( n \)-types. Let \( \phi_i \) be a generating formula for \( X_i \). Then, if \( i \neq j \) we must have that \( \neg \phi_i \in X_j \). In particular, if \( j > \max\{i_1, \ldots, i_k\} \) we will have that \( \neg(\phi_{i_1} \land \ldots \land \phi_{i_k}) \in X_j \).

It follows that \( \{ \neg \phi_i ; i \in \mathbb{N} \} \) is a consistent set (here we use the assumption that there are infinitely many principal \( n \)-types), and may thus be extended to an \( n \)-type different from all the \( X_i \)'s. This contradicts the assumption that we have enumerated them all, and the lemma is proved.

If \( \mathfrak{A} \) is a model for \( T \) and \( a_1, \ldots, a_n \) are elements in \( A \) (possibly with repetition), then \( a_1, \ldots, a_n \) realizes exactly one \( n \)-type,
\[
X = \{ \phi(x_1, \ldots, x_n) ; \mathfrak{A} \models \phi(a_1, \ldots, a_n) \}.
\]
This indicates that the set of types may give us information about the class of countable structures. One example of this is

**Theorem 1.4.8** Assume that all \( n \)-types are principal for all \( n \).

Then \( T \) is \( \omega \)-categorical.

**Proof**
We will elaborate on the proof of the \( \omega \)-categoricity for the theory \( DO \).

Let \( \mathfrak{A} \) and \( \mathfrak{B} \) be two countable models for \( T \) with domains \( A \) and \( B \). Let \( A' = \{ a_1, \ldots, a_n \} \subseteq A \) and \( B' = \{ b_1, \ldots, b_n \} \subseteq B \) be finite sets (without repetition) and let \( p(a_i) = b_i \). We call \( p \) a partial isomorphism if \( \{ a_1, \ldots, a_n \} \) and \( \{ b_1, \ldots, b_n \} \) realize the same \( n \)-type, we have that

**Claim**
If \( p \) is a partial isomorphism and \( a \in A \), then \( p \) can be extended to a partial isomorphism \( q \) defined on \( a \).

Let us first see how to prove the theorem from the claim. The statement is symmetrical in \( \mathfrak{A} \) and \( \mathfrak{B} \), so given \( b \in B \) we may as well extend \( p \) to \( a \) with \( b \) in its range. \( \mathfrak{A} \) and \( \mathfrak{B} \) will be elementary equivalent since \( T \) is complete. We may consider the empty function as a partial isomorphism if \( \{ a_1, \ldots, a_n \} \) and \( \{ b_1, \ldots, b_n \} \) realize the same \( n \)-type, we have that
each $a \in A$ will be in the domain of one of the partial isomorphisms, and each $b \in B$ will be in the range of one of the partial isomorphisms. In the end, we have constructed a total isomorphism.

**Proof of claim**

Let $X$ be the $n+1$-type of $\{a_1, \ldots, a_n, a\}$. Since we have assumed that all $n+1$-types are principal, there is a generating formula $\phi(x_1, \ldots, x_{n+1})$ in $X$.

Since $\mathfrak{A} \models \phi(a_1, \ldots, a_n, a)$ we have that
\[
\mathfrak{A} \models \exists x_{n+1} \phi(a_1, \ldots, a_n, x_{n+1}).
\]

Since $\{a_1, \ldots, a_n\}$ and $\{b_1, \ldots, b_n\}$ realize the same $n$-type, we have that
\[
\mathfrak{B} \models \exists x_{n+1} \phi(b_1, \ldots, b_n, x_{n+1}).
\]

Choose $b \in B$ such that $\phi(b_1, \ldots, b_n, b)$ holds in $\mathfrak{B}$ and let $q(a) = b$.

Let $\psi \in X$. Since $T \vdash \phi(x_1, \ldots, x_n, x_{n+1}) \rightarrow \psi(x_1, \ldots, x_n, x_{n+1})$ we see that $\psi(b_1, \ldots, b_n, b)$ will hold in $\mathfrak{B}$. This shows that $X$ is the type realized by $\{b_1, \ldots, b_n, b\}$, and the claim is proved.

As a consequence we see that if the number of $n$-types is finite for each $n$, then $T$ is $\omega$-categorical. We will prove the converse to this, but in order to do so we need the theorem about omitting types. The proof of the omitting type theorem is based on the Henkin-style proof of the completeness theorem for countable, first order theories. In that proof we we first extend a theory $T$ to a Henkin theory $T_H$ by adding Henkin constants $c_1, c_2, \ldots$ and Henkin axioms
\[
\exists x \phi \rightarrow \phi_{c_1}^x,
\]
and then extend $T_H$ to a complete theory $T_C$. The model realizing the completeness theorem is then the term model of $T_C$. In this argument there is an amount of freedom, since there may be many ways to construct $T_C$ from $T_T$. This freedom will be used in the proof of the omitting type theorem.

**Theorem 1.4.9** Let $X$ be a non-principal $n$-type. Then there is a model for $T$ omitting $X$.

**Proof**

The idea of the proof is as follows: We extend the language $L$ of $T$ with Henkin constants $c_1, c_2, \ldots$ and we extend $T$ with Henkin axioms. Then, following the proof of the completeness theorem, we must make a complete extension of this extended theory and then form the term model. The term model will consist of equivalence classes of closed terms, and in fact, via the Henkin constant for $\exists x (x = t)$, we see that each equivalence class will contain a Henkin constant. Thus we must ensure, during the completion process, that for any
collection $c_1, \ldots, c_n$ of Henkin constants there is a formula $\phi \in X$ such that
$\neg \phi(c_1, \ldots, c_n)$ is added to the theory.

In the proof of the completeness theorem we add Henkin constants and Henkin axioms in waves. However, we prove that each wave is countable. Thus the set of new constants is countable, and we may just organize them in a list as above such that if $c_j$ is the Henkin constant of $\exists x \phi$ and $c_i$ occurs in $\phi$, then $i < j$.

Recall that the Henkin axiom is
$\exists x \phi(x) \rightarrow \phi(c_j)$.

For the sake of simplicity, we will assume that $n = 1$.
For each Henkin constant $c_k$ we intend to find a formula $\phi_k \in X$ such that
$T \cup \{\neg \phi_1(c_1), \ldots, \neg \phi_k(c_k)\}$
is consistent with all Henkin axioms.
From Henkin’s proof of the completeness theorem we know that adding a Henkin axiom introducing a new Henkin constant will preserve consistency. Thus, in the process of verifying that $\phi_k(c_k)$ is consistent with the extended theory so far and all the Henkin axioms, we only need to be concerned with the Henkin axioms involving $c_1, \ldots, c_k$.

The ‘construction’ is by recursion on $k$, so assume that
$T \cup \{\neg \phi_1(c_1), \ldots, \neg \phi_{k-1}(c_{k-1})\}$
is consistent with all the Henkin axioms.
Let $\xi(c_1, \ldots, c_k)$ be the conjunction of all $\neg \phi_j(c_j)$ for $j < i$ and all Henkin axioms introducing $c_1, \ldots, c_k$.
Then $T, \xi(c_1, \ldots, c_k)$ is consistent.
With the aim of obtaining a contradiction, assume that $T, \xi(c_1, \ldots, c_k), \neg \phi(c_k)$ is inconsistent for all $\phi \in X$.
Then $T \vdash \xi(c_1, \ldots, c_k) \rightarrow \phi(c_k)$ for all $\phi \in X$.

By the theorem of constants
$T \vdash \xi(y_1, \ldots, y_{k-1}, x) \rightarrow \phi(x)$
and thus
$T \vdash \exists y_1 \cdots \exists y_{k-1} \xi(y_1, \ldots, y_{k-1}, x) \rightarrow \phi(x)$
for all $\phi \in X$. Since $X$ is complete, there are two possibilities

1. $\exists y_1 \cdots \exists y_{k-1} \xi(y_1, \ldots, y_{k-1}, x) \in X$

2. $\neg \exists y_1 \cdots \exists y_{k-1} \xi(y_1, \ldots, y_{k-1}, x) \in X$.

Case 1. contradicts that $X$ is non-principal.
In case 2., we let $\phi = \neg \exists y_1 \cdots \exists y_{k-1} \xi(y_1, \ldots, y_{k-1}, x)$ and see that

$T \vdash \exists y_1 \cdots \exists y_{k-1} \xi(y_1, \ldots, y_{k-1}, x) \rightarrow \neg \exists y_1 \cdots \exists y_{k-1} \xi(y_1, \ldots, y_{k-1}, x)$.
By propositional logic it follows that

\[ T \vdash \neg \exists y_1 \cdots \exists y_{k-1} \xi(y_1, \ldots, y_{k-1}, x) \]

contradicting the induction hypothesis, i.e. that \( T, \xi(c_1, \ldots, c_k) \) is consistent.

The assumption was that \( T, \xi(c_1, \ldots, c_k), \neg \phi(c_k) \) is inconsistent for each \( \phi \in X \). This led to a contradiction, so there is a \( \phi_k \in X \) such that

\[ T, \xi(c_1, \ldots, c_k), \neg \phi_k(c_k) \]

is consistent.

The choice of \( \xi \) was such that \( T, \neg \phi_1(c_1), \ldots, \neg \phi_k(c_k) \) will be consistent with all the Henkin axioms.

Let \( T^\omega \) be \( T \) extended with all \( \neg \phi_k(c_k) \) and all Henkin axioms. \( T^\omega \) will be a consistent Henkin theory. Let \( T^c \) be a completion of \( T^\omega \). Then the term model of \( T^c \) will be a model for \( T \) omitting \( X \).

In Exercise 1.24 we will discuss why the assumption that \( n = 1 \) is a harmless one.

**Corollary 1.4.10** If there are infinitely many \( n \)-types for some \( n \), then \( T \) is not \( \omega \)-categorical.

**Proof**

If there are infinitely many \( n \)-types, one of them must be non-principal. Let \( X \) be non-principal. Then there is one model \( \mathfrak{A} \) realizing \( X \) and one model \( \mathfrak{B} \) omitting \( X \). These models cannot be isomorphic.

**Supplementary material for the advanced reader**

This tiny section can only be read with complete understanding by a reader with a background from descriptive set theory and cardinality arithmetics.

An \( n \)-type is a subset of the set of formulas \( \phi(x_1, \ldots, x_n) \), which is a countable set. Via e.g. a Gödel enumeration we may consider an \( n \)-type to be a subset of the natural numbers. The set of subsets of \( \mathbb{N} \) forms a compact, metrizable topological space in a natural way, homeomorphic to the Cantor space.

There are two requirements that have to be satisfied by an \( n \)-type. Given

\[ \phi_1, \ldots, \phi_k \in X \]

we must have that \( T \vdash \exists x_1 \cdots \exists x_n (\phi_1 \land \cdots \land \phi_k) \), and \( X \) must be complete. The first requirement defines a \( G_\delta \)-set and the second a closed set. Thus the set of \( n \)-types is a \( G_\delta \)-set of sets of formulas.

By a standard and classical result of descriptive set theory, any \( G_\delta \)-set is either finite, countable or has the cardinality of the continuum (this holds for a much more general class than the \( G_\delta \)-sets). Thus, if there are uncountably many \( n \)-types for some \( n \) there will be as many isomorphism classes of countable models for \( T \) as there are reals, since \( \kappa \) many countable models will at most realize \( \omega \times \kappa \) many \( n \)-types, and if \( \kappa < 2^\omega \) then \( \omega \times \kappa < 2^\omega \).
1.5 Saturated structures

In a sense we may say that the more \( n \)-types a model \( \mathfrak{A} \) for \( T \) realize, the richer the model is. The most generous would have been if \( \mathfrak{A} \) realizes every \( n \)-type for every \( n \). Of course, if there are uncountably many \( n \)-types for some \( n \), this is impossible, a countable structure can realize at most countably many \( n \)-types.

**Definition 1.5.1** Let \( \mathfrak{A} \) be a model for \( T \), \( A_0 \subseteq A \).

Let \( L(A_0) \) be \( L \) extended with names for each \( a \in A_0 \). By abuse of notation we will use \( a \) both as the name and for the object. Let \( T[A_0] \) be the complete theory of all \( L(A_0) \)-sentences that are true in \( \mathfrak{A} \).

There is no reference to \( \mathfrak{A} \) in the notation ‘\( T[A_0] \)’. This will be harmless in the proof of Theorem 1.5.6 by the following observation:

**Lemma 1.5.2** Let \( \mathfrak{A} \) be an elementary substructure of \( \mathfrak{B} \) and let \( A_0 \subseteq A \) be finite.

Then \( T[A_0] \) will be the same if we view \( A_0 \) as a subset of \( A \) or as a subset of \( B \).

**Proof**
This Lemma is just a reformulation of the definition of elementary substructures.

**Definition 1.5.3** Let \( \mathfrak{A} \) be a model for \( T \). We call \( \mathfrak{A} \) saturated if all \( 1 \)-types in \( T[A_0] \) are realized in \( \mathfrak{A} \) whenever \( A_0 \) is finite.

**Lemma 1.5.4** Assume that \( T \) has only countably many \( n \)-types for each \( n \). Let \( \mathfrak{A} \models T \). Let \( A_0 \subseteq A \) be finite.

Then \( T[A_0] \) has only countably many \( 1 \)-types.

**Proof**
Let \( A_0 = \{a_1, \ldots, a_n\} \) and let \( X \) be a \( 1 \)-type in \( T[A_0] \).

Let \( X' = \{\phi(x_1, \ldots, x_n, x) : \phi(a_1, \ldots, a_n, x) \in X\} \).

Clearly, \( X' \) is a consistent set, so \( X' \) can be extended to an \( n+1 \)-type.

If \( X \) and \( Y \) are different \( 1 \)-types in \( T[A_0] \), then for some \( \phi(x) \) in \( L(A_0) \), \( \phi \in X \land \neg \phi \in Y \). It follows that \( X' \) and \( Y' \) cannot be extended to the same \( n+1 \)-type. Since there are at most countably many \( n+1 \)-types in \( T \), there are at most countably many \( 1 \)-types in \( T[A_0] \).

**Lemma 1.5.5** Let \( \mathfrak{A} \) be a model for \( T \), \( A_0 \subseteq A \) a finite set and \( X \) a \( 1 \)-type in \( T[A_0] \). Then there is an elementary extension of \( \mathfrak{A} \) realizing \( X \).

**Proof**
If \( X \) is principal, \( \mathfrak{A} \) will realize \( X \). So assume that \( X \) is not principal. Then \( \mathfrak{A} \) must be infinite (why?). Extend \( L \) by constants \( c_a \) for each element \( a \in A \), and extend \( T \) to the complete theory \( Th(\mathfrak{A}) \).
Let \( c \) be a new constant and extend \( Th(\mathfrak{A}) \) to \( T^* \) by adding the axioms

\[ \phi(c) \text{ whenever } \phi \in X. \]
\[ c \neq c_a \text{ for all } a \in A. \]

We then use the compactness theorem to see that \( T^* \) has a model \( \mathfrak{B} \). \( \mathfrak{B} \) will be an elementary extension of \( \mathfrak{A} \) realizing \( X \).

**Theorem 1.5.6** Assume that \( T \) has at most countably many \( n \)-types for each \( n \in \mathbb{N} \).

Then \( T \) has a saturated model.

**Proof**

The final model will be the directed limit of an elementary directed system, and if the reader has not solved Exercise 1.7 so far it is about time to do so before any further reading of this proof.

We will start with a countable model \( \mathfrak{A}_0 \) with domain \( A_0 \). At stage \( k \), assume that we have constructed a countable model \( \mathfrak{A}_k \) with domain \( A_k \). Then we select a finite subset \( B_k \subseteq A_k \) and a 1-type \( X_k \) in \( T[B_k] \) and let \( \mathfrak{A}_{k+1} \) be an elementary extension of \( \mathfrak{A}_k \) realizing \( X_k \). Here we use Lemma 1.5.5.

In order to obtain a saturated model in the end, we must organize the selection of \( B_k \) and \( X_k \) such that the following are satisfied:

- Whenever \( B \subseteq A_k \) is finite, there are infinitely many \( k' \geq k \) such that \( B = B_{k'} \).
- Whenever \( B \subseteq A_k \) is finite and \( X \) is a 1-type in \( T[B] \), then there is a \( k' \geq k \) such that \( B = B_{k'} \) and \( X = X_{k'} \).

This organization requires some standard techniques, and the reader acquainted with those may skip the rest of the proof. For the rest of the readers, we offer the details:

When \( \mathfrak{A}_k \) is constructed, let \( \{C_{i,k}\}_{i \in \mathbb{N}} \) be an enumeration of all the finite subsets of \( A_k \) (that are not subsets of \( A_{k-1} \) when \( k > 0 \); this is not essential). By Lemma 1.5.2 \( T[C_{i,k}] \) will be a fixed theory for the rest of the construction.

For each \( i \) and \( k \), let \( \{X_{j,i,k}\}_{j \in \mathbb{N}} \) be an enumeration of all 1-types in \( T[C_{i,k}] \).

If \( k' \geq 1 \), we may write \( k' \) in a unique way as

\[ k' = 2^a3^b5^c m \]

where \( m \) is not divisible by 2, 3 or 5. Then we let \( B_{k'} = C_{i,k} \) and \( X_{k'} = X_{j,i,k} \).

**Theorem 1.5.7** Let \( \mathfrak{A} \) and \( \mathfrak{B} \) be countable models for \( T \) where \( \mathfrak{B} \) is saturated.

Then there is an elementary embedding from \( \mathfrak{A} \) to \( \mathfrak{B} \).

If both \( \mathfrak{A} \) and \( \mathfrak{B} \) are saturated, they are isomorphic.

**Proof**

We show how to construct an elementary embedding in the first case. The isomorphism in the second case can be constructed by a back-and-forth construction using our construction at each step.

Let \( A \) and \( B \) be the domains of \( \mathfrak{A} \) and \( \mathfrak{B} \). Let \( A = \{a_k : k \in \mathbb{N}\} \), let
$A_k = \{a_0, \ldots, a_{k-1}\}$ and let $\pi_k : A_k \to B$ be injective.
We say that $\pi_k$ is a partial elementary embedding if for all sentences $\phi$ in $L(A_k)$,
\[
\mathfrak{A} \models \phi \iff \mathfrak{B} \models \phi^{\pi_k}
\]
where $\phi^{\pi_k}$ is the sentence in $L(\mathfrak{B})$ obtained from $\phi$ by replacing each name for an element $a \in A_k$ by the name for $\pi_k(a)$. Observe that for $k = 0$, $\pi_k$ is the empty function, and then it will be a partial elementary embedding just because $\mathfrak{A}$ and $\mathfrak{B}$ are elementary equivalent.

Now let $B_k = \{\pi_k(a) : a \in A_k\}$.
Let $X$ be the 1-type in $T[A_k]$ realized by $a_k$, and let $X'$ be the corresponding 1-type in $T[B_k]$ using the transformation $\phi \mapsto \phi^{\pi_k}$.
Since $\mathfrak{B}$ is saturated, $X'$ is realized by some $b \in B$.
We extend $\pi_k$ to $\pi_{k+1}$ by $\pi_{k+1}(a_k) = b$ where $b \in B$ realizes $X'$. Then $\pi_{k+1}$ is a partial elementary embedding and the construction can go on.
At the end, the union $\pi$ of all the $\pi_k$’s will be an elementary embedding from $\mathfrak{A}$ to $\mathfrak{B}$.

The number of countable models

The possible cardinalities of the set of isomorphism classes of countable models of a complete countable theory is fully determined, the cardinalities can be

$1, 3, 4, \ldots, \aleph_0, \aleph_1, 2^{\aleph_0}$.

That the number 2 is left out is not a typo, but is one of the more peculiar results in model theory, a result we will discuss further.
The theory $DO^+$ has three different models. Elaborating on this kind of construction we may find theories where the number of models are any number $\neq 2$.
The theory $ACF(0)$ has countably many models. There are examples of theories with continuumly many models. One example is $DO !$, an extension of $DO$ defined as follows:

- Let $\{c_{i,j}\}_{i,j \in \mathbb{N}}$ be new constants.
- Let $c_{i,j} < c_{i',j'}$ be an axiom if $i < i'$ or if $i = i'$ and $j < j'$.

Then $DO !$ is a complete theory because the restriction of $DO !$ to a finite set of constants will be $\omega$-saturated.
$DO !$ will have continuumly many non-isomorphic countable models. Each model $\mathfrak{A}$ will determine a subset $X_{\mathfrak{A}}$ of $\mathbb{N}$ by $n \in X_{\mathfrak{A}}$ if there is an element between all the $c_{n,j}$’s and $c_{n+1,0}$. All sets can be realized this way, and isomorphic structures clearly define the same set.

In order to prove that other cardinalities are out of the question, we need to go beyond first order logic, and we need background theorems from descriptive set theory. This is outside the scope of this compendium. It is just recently established (2002) that $\aleph_1$ is a possible alternative, even in the case that the
continuum hypothesis fails. The construction behind this argument is far beyond our ambitions here.
We do however have the tools available in order to show that there are never exactly 2 isomorphism classes of models. We give a sketch of the proof. The reader is invited to fill in the details in the (minor) gaps in the proof of Lemma 1.5.9 and to prove Corollary 1.5.10.

**Definition 1.5.8** A model $\mathfrak{A}$ for $T$ is *weakly saturated* if each $n$-type is realized in $\mathfrak{A}$ for each $n \in \mathbb{N}$.

**Lemma 1.5.9** Assume that the number of models for $T$ is finite, but $> 1$. Then $T$ has a model that is weakly saturated, but not saturated.

**Sketch of proof**
$T$ has a countable, saturated model $\mathfrak{A}$. Let $\mathfrak{A}_1, \ldots, \mathfrak{A}_k$ be the other countable models. Assume that neither of these are weakly saturated, and for each $i \leq k$, let $X_i$ be an $n_i$-type not realized in $\mathfrak{A}_i$. Without loss of generality, we may assume that the variables used for the types $X_i$ are distinct, so $X_1 \cup \cdots \cup X_k$ can be viewed as a set of formulas in the variables e.g. $x_1, \ldots, x_n$ where $n = \sum_{i=1}^k n_i$. Let $X$ be an $n$-type extending $X_1 \cup \cdots \cup X_k$. Then $X$ cannot be realized in $\mathfrak{A}_i$ for any $i \leq k$.
Let $c_1, \ldots, c_n$ be new names, and let $T^*$ be $T$ extended with the axioms $\phi(c_1, \ldots, c_n)$ for $\phi(x_1, \ldots, x_n) \in X$. Since $\mathfrak{A}$ is saturated, we may interpret the constants $c_1, \ldots, c_n$ in $\mathfrak{A}$ in such a way that $\mathfrak{A}$ becomes a model for $T^*$. Moreover, two such models will be isomorphic. Thus $T^*$ is $\omega$-categorical and there are finitely many $n$-types in $T^*$ for each $n$. It follows that there are only finitely many $n$-types in $\mathfrak{T}$ for each $n$, which contradicts the assumption.
This ends the sketch.

**Corollary 1.5.10** There is no complete countable first order theory with exactly two isomorphism classes of countable models.

This corollary follows from the lemma and most of the results established earlier in this section, but the proof will not require any new technical arguments. There is no complete countable theory with exactly two non-isomorphic countable models.

### 1.6 $\omega$-logic

The prime object of consideration in this chapter has been the model theory of first order logic. We have seen that for many purposes, first order logic is inadequate. Logicians have been studying other forms of logic where some of the tools from first order logic can be used. There are at least three directions to go, we may extend the use of quantifiers beyond the setting of first order logic, we may consider infinitary proof trees and we may consider infinite connectives. The first alternative is natural if we want to form logics adequate for mathematical analysis or e.g. for the study of Noetherian rings. The third alternative
was used when the possible numbers of models for a complete first order theory was characterized in full, something we chose not to do.

In this section we will consider a form of infinitary logic for models where the natural numbers is a definable substructure. This actually requires infinite proof trees, since we cannot have the compactness theorem to be valid in this case.

1.6.1 \(\omega\)-logic

For the rest of this section, let \(L\) be a first order language with (among possibly other symbols) a unary predicate symbol \(\Omega\), a constant 0 and a unary function symbol \(S\). We will let \(T\) be a first order theory over the language \(L\).

**Definition 1.6.1** Let \(\mathfrak{A}\) be a model for \(T\).

We call \(\mathfrak{A}\) an \(\omega\)-model if \(\Omega\) is interpreted as \(\mathbb{N}\), 0 as zero and \(S\) as the successor-function.

**Definition 1.6.2** Let \(T\) be a theory as above. \(T\) is \(\omega\)-sound if \(T\) has an \(\omega\)-model.

A formula \(\phi\) in \(L\) is \(\omega\)-valid if \(\phi\) is valid in all \(\omega\)-models for \(T\).

We may formulate a logical system adequate for this concept of validity:

**Definition 1.6.3** \(\omega\)-logic will be first order logic extended with four axioms and one infinitary rule:

- Let \(k_0\) be the constant 0, and let \(k_{n+1}\) be the term \(S(k_n)\).
- \(\Omega(0)\) and \(\Omega(x) \rightarrow \Omega(S(x))\) are new axioms in \(\omega\)-logic.
- \(0 \neq S(x)\) and \(S(x) = S(y) \rightarrow x = y\) are new axioms in \(\omega\)-logic.
- If \(\phi^x_n\) is an \(\omega\)-theorem in \(T\) for all \(n \in \mathbb{N}\), then \(\forall x(\Omega(x) \rightarrow \phi)\) is an \(\omega\)-theorem in \(T\).
- We will write \(T \vdash_\omega \phi\) when \(\phi\) is an \(\omega\)-theorem in \(T\).

**Remark 1.6.4** We will not give a precise definition of an \(\omega\)-proof, because in order to do so we need to develop a theory of infinite, well founded trees. The important fact is that we have used recursion to define the class of \(\omega\)-theorems, and then we can prove lemmas and theorems about them using induction over this construction.

In Exercise 5.14 we will make use of the concept of well founded trees introduced in Chapter 4 and give a precise definition of an \(\omega\)-proof. The present section will only be technically sound on the basis of this exercise.

\(\omega\)-logic is of interest for several reasons. One reason is that the induction axiom in first order number theory is provable in \(\omega\)-logic. Assume that we have a proof of

\[
\phi(0)
\]
and that we have a proof of
\[ \forall x (\Omega(x) \land \phi(x) \rightarrow \phi(S(x))). \]
We then clearly have \( \omega \)-proofs for each \( \phi(k_n) \) and by the \( \omega \)-rule we have a proof of
\[ \forall x (\Omega(x) \rightarrow \phi(x)). \]
In some odd respect, this is a simplification; the \( \omega \)-logic version of Peano arithmetic is better behaved from a proof-theoretical point of view, and analyzing the \( \omega \)-logic versions of proofs in number theory gives us information about the strength of number theory. This kind of application requires a much deeper involvement in proof theory than we are prepared for in this compendium.

### 1.6.2 \( \omega \)-completeness

We will prove the completeness theorem for \( \omega \)-logic. One motivation is to demonstrate the power of the proof method of the first order version. We will leave many of the details as exercises for the reader.

**Definition 1.6.5** \( T \) is \( \omega \)-consistent if there is no sentence \( \phi \land \neg \phi \) that is an \( \omega \)-theorem in \( T \).

We have the standard equivalent versions of the completeness theorem for \( \omega \)-logic:

**Lemma 1.6.6** The following are equivalent:

1. A formula \( \phi \) is valid in all \( \omega \)-models for \( T \) if and only if \( \phi \) is an \( \omega \)-theorem in \( \phi \).

2. A theory \( T \) is \( \omega \)-consistent if and only if \( T \) has an \( \omega \)-model.

The proof is left as Exercise 1.25.

**Lemma 1.6.7** The theorem of constants holds for \( \omega \)-logic.

The proof is left as Exercise 1.26.

**Lemma 1.6.8** Let \( L \) be as above, \( T \) be an \( \omega \)-consistent theory over \( L \) and \( \exists x \phi \) be a sentence in \( L \). Let \( r \) be a new name, not in \( L \). Then \( T, \exists x \phi \rightarrow \phi^r \) is also \( \omega \)-consistent.

The proof is left as Exercise 1.27.

Using this lemma, we may construct the Henkin-extension of an \( \omega \)-consistent theory such that all finite subtheories are \( \omega \)-consistent as well. We cannot conclude that the full Henkin extension is \( \omega \)-consistent, simply since \( \omega \)-consistency is not closed under directed unions of theories. We will see how we can avoid this obstacle.
From now on, \( T \) is a fixed \( \omega \)-consistent theory. Let \( r_0, r_1, \ldots \) be the new constants that we need in order to construct the Henkin extension \( H \) of \( T \), and let \( L_\infty \) be the extended language. Let \( H_n \) be the subtheory of \( H \) where we only added the constants \( r_0, \ldots, r_{n-1} \) to \( T \) and the corresponding Henkin axioms.

Then \( r_n \) is not in the language of \( H_n \).

Let \( \{ \phi_n \}_{n \in \mathbb{N}} \) be an enumeration of all sentences in \( L_\infty \) such that \( \phi_n \) is a sentence in the language of \( H_n \).

**Lemma 1.6.9** Let \( S \) be an \( \omega \)-consistent theory, \( \phi \) a sentence in the language of \( S \). Then at least one of the theories \( S, \phi \) and \( S, \neg \phi \) are \( \omega \)-consistent.

The proof is left as Exercise 1.28.

**Lemma 1.6.10** Let \( S \) be an \( \omega \)-consistent theory, \( t \) a closed term in the language of \( S \) and assume that \( S, \Omega(t) \) is \( \omega \)-consistent.

Then there is a number \( n \) such that \( S, t = k_n \) is \( \omega \)-consistent.

**Proof**

Assume that \( S, t = k_n \) is inconsistent for all \( n \).

Then \( t \neq k_n \) is an \( \omega \)-theorem in \( S \) for all \( n \).

Using the \( \omega \)-rule we obtain

\[ S \vdash_\omega \forall x(\Omega(x) \rightarrow x \neq t). \]

This means that \( S \vdash_\omega \neg \Omega(t) \), contradicting the assumption. This ends the proof of the lemma.

We will now construct the theory \( T_\infty \).

Let \( T_0 = H_0 = T \).

Assume that \( T_n \) is an \( \omega \)-consistent extension of \( H_n \) within the language of \( H_n \).

By Lemma 1.6.8, \( T \cup H_{n+1} \) will be \( \omega \)-consistent.

If \( T_n \cup H_{n+1}, \phi_n \) is not \( \omega \)-consistent, we let

\[ T_{n+1} = T_n \cup H_{n+1}, \neg \phi_n. \]

If \( T_n \cup H_{n+1}, \phi_n \) is \( \omega \)-consistent, but \( \phi_n \) is not of the form \( \Omega(t) \), let

\[ T_{n+1} = T_n, H_{n+1}, \phi_n. \]

If \( T_n \cap H_{n+1}, \phi_n \) is \( \omega \)-consistent, and \( \phi_n \) is of the form \( \Omega(t) \), we choose one \( m \) such that

\[ T_{n+1} = T_n \cup H_{n+1}, t = k_m \]

is \( \omega \)-consistent.

Let \( T_\infty \) be the union of all \( T_n \). We cannot prove directly that \( T_\infty \) is \( \omega \)-consistent, but \( T_\omega \) is a complete Henkin theory in the traditional sense, so the term model \( \mathcal{A} \) of \( T_\infty \) will be a model for \( T_\infty \), and then in particular for \( T \).

\( \Omega^A \) will be the set of equivalence classes of closed terms \( t \) such that \( T_\infty \vdash \Omega(t) \).

For each such \( t \), there will be a number \( n \in \mathbb{N} \) such that \( T_\infty \vdash t = k_n \). Thus \( \Omega^A \) will be isomorphic to \( \mathcal{A} \). This indirect argument shows that \( \mathcal{A} \) is an \( \omega \)-model after all, and the completeness theorem for \( \omega \)-logic is proved. Thus, \( T_\infty \) is actually \( \omega \)-consistent.
1.7 Exercises to Chapter 1

Exercise 1.1 Show that if \( \pi \) is an isomorphism, then \( \pi \) has an inverse \( \pi^{-1} \) that is also an isomorphism.

Exercise 1.2 Show that if \( T \) is an open theory, \( \mathfrak{A} \) is a model for \( T \) and \( \pi : \mathfrak{B} \rightarrow \mathfrak{A} \) is an embedding, then \( \mathfrak{B} \) is a model for \( T \). (This is the easy part of Los-Tarski’s theorem.)

Exercise 1.3 Let \( L \) be a first order language, \( \mathfrak{A} \) an \( L \)-structure and \( \phi \) an \( L \)-formula with at most \( x_1, \ldots, x_n \) free.

a) Let \( s \) be an assignment over \( \mathfrak{A} \). Then
\[
\mathfrak{A} \models t^\mathfrak{A}[s] = \left( t_{x_1, \ldots, x_n}^{C_{x_1}(x_1), \ldots, C_{x_n}(x_n)} \right)^\mathfrak{A}
\]
whenever \( t \) is a term with at most \( x_1, \ldots, x_n \) free.

\[
\mathfrak{A} \models \phi[s] \Leftrightarrow \mathfrak{A} \models \phi_{x_1, \ldots, x_n}^{C_{x_1}(x_1), \ldots, C_{x_n}(x_n)}
\]
whenever \( \phi \) is a formula with at most \( x_1, \ldots, x_n \) free.

This shows that the definition of interpretations via assignments and via new names for all elements of the structure in question are equivalent.

b) Let \( t \) be a term with variables among \( x_1, \ldots, x_n \), and let \( u_1, \ldots, u_n \) and \( r_1, \ldots, r_n \) be two sequences of closed terms such that \( u_i^\mathfrak{A} = r_i^\mathfrak{A} \) for each \( i \).

Then \( (t_{u_1, \ldots, u_n})^\mathfrak{A} = (t_{r_1, \ldots, r_n})^\mathfrak{A} \).

c) Let \( \phi \) be a formula with free variables among \( x_1, \ldots, x_n \), and let \( u_1, \ldots, u_n \) and \( r_1, \ldots, r_n \) be two sequences of closed terms such that \( u_i^\mathfrak{A} = r_i^\mathfrak{A} \) for each \( i \).

Then \( (\phi_{u_1, \ldots, u_n})^\mathfrak{A} \Leftrightarrow (\phi_{r_1, \ldots, r_n})^\mathfrak{A} \).

Exercise 1.4 Let \( \mathfrak{A} \) and \( \mathfrak{B} \) be given.
Show that the following are equivalent:

1. There is an embedding \( \pi \) from \( \mathfrak{A} \) to \( \mathfrak{B} \).
2. \( \mathfrak{A} \) is isomorphic to a substructure \( \mathfrak{A}' \) of \( \mathfrak{B} \).
3. \( \mathfrak{A} \) is a substructure of some \( \mathfrak{B}' \) isomorphic to \( \mathfrak{B} \).

Exercise 1.5 Work out a complete proof of Theorem 1.1.14. In particular:

a) Show that if both \( \mathfrak{A} \) and \( \mathfrak{B} \) satisfy the requirements of a directed limit, then they are isomorphic.

b) Show that the definition of \( \equiv \) is independent of the choice of \( k \).

c) Show that if \( C_1, \ldots, C_n \) are equivalence classes, there is a common \( i \in I \) such that each \( C_j \) have an element of the form \( (i, a_j) \).
d) Show that the definitions of $f^A$ and $R^A$ are independent of the choice of $i$.

e) Show that $\eta$ as constructed in the proof is an embedding satisfying the requirement of the theorem, and that this is the only one doing so.

Then write down a full proof where a) - e) are incorporated in a natural way.

**Exercise 1.6** Let $L$ be a first order language (with equality) and let $T$ be a complete theory over $L$.

a) Show that either is each model of $T$ infinite or all models of $T$ will have the same finite cardinality.

b) Show that if $T$ has a finite model, then all models are isomorphic.

Hint: Use Theorem 1.2.11.

**Exercise 1.7** Let $\langle \{\mathfrak{A}_i\}_{i \in I}, \{\pi_{ij}\}_{i<j}\rangle$ be an elementary directed system with directed limit $\langle \mathfrak{A}, \{\pi_i\}_{i \in I}\rangle$.

Uniformly in $i \in I$, show by induction on the complexity of a sentence $\phi$ in $L(\mathfrak{A}_i)$ that $\pi_i$ preserves the truth value of $\phi$.

**Exercise 1.8** Show that if $T$ is a countable, $\omega$-categorical theory, then $T$ is complete.

**Exercise 1.9** Prove in detail that DO is $\omega$-categorical

**Exercise 1.10**

a) Show that for each prime number $p$ there is a finite extension $F(p)$ of field theory such that the models are exactly the fields of characteristic $p$.

b) Show that there is an extension $F(0)$ of field theory such that the models are exactly the fields of characteristic 0.

c) Let $\phi$ be a sentence in field theory. Show that if $\phi$ is true in fields of arbitrarily large finite characteristic, then $\phi$ is valid in some field of characteristic 0.

d) Let $\phi$ be a formula in field theory. Show that if $\phi$ is valid in finite fields of arbitrarily large cardinality, then $\phi$ is valid in some infinite field.

e) Show that the theory of finite fields or the theory of fields of finite characteristics cannot be expressed in first order logic.

f) Show that the theory $F(0)$ or the theory of infinite fields are not finitely axiomatizable.

**Exercise 1.11** (For readers familiar with Zorn’s Lemma).

Let $\mathfrak{A}$ and $\mathfrak{B}$ be algebraically closed fields and let $\mathfrak{A}_0$ and $\mathfrak{B}_0$ be isomorphic substructures via the isomorphism $\pi_0$. 

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a) Use Zorn's lemma to show that $\pi_0$ can be extended to a maximal isomorphism $\pi_1$ between substructures $\mathfrak{A}_1$ and $\mathfrak{B}_1$.

b) Show that $\mathfrak{A}_1$ and $\mathfrak{B}_1$ are algebraically closed subfields of $\mathfrak{A}$ and $\mathfrak{B}$.

**Exercise 1.12** Let $DO^+$ be the theory in Example 1.3.22.

a) The theory $DO^+$ has three countable, non-isomorphic models. Find examples of these, and show that there are no more than three isomorphism classes of countable models.

Hint: When you have found three non-isomorphic examples, you may use the proof of the $\omega$-categoricity of $DO$ to prove that there are no more.

b) Does $DO^+$ have the isomorphism property?

c) Show that $DO^+$ is complete.

Hint: Show that certain finite subtheories of $DO^+$ are $\omega$-categorical.

**Exercise 1.13**

a) Let $T$ be a first order theory with both the isomorphism and submodel properties. Let $\mathfrak{A}$ be a model for $T$.

Show that $T(\mathfrak{A}) = T \cup D(\mathfrak{A})$ has the isomorphism and submodel properties.

b) A theory $T$ is called *model complete* if $T \cup D(\mathfrak{A})$ is complete whenever $\mathfrak{A} \models T$.

Show that if $T$ is consistent and has the isomorphism and submodel properties, then $T$ is model complete.

**Exercise 1.14** Let $\mathfrak{F}$ be a field, $P_i(x_1, \ldots, x_k)$ and $Q_j(x_1, \ldots, x_k)$ be polynomials over $\mathfrak{F}$ for $i \leq n$ and $j \leq m$.

Show that there is a solution to the simultaneous set of equations

$$P_i(x_1, \ldots, x_k) = 0$$

and inequations

$$Q_j(x_1, \ldots, x_k) \neq 0$$

in some field extending $\mathfrak{F}$ if and only if there is a solution in the algebraic closure of $\mathfrak{F}$.

(This result is known as Hilbert’s Nullstellensatz.)

**Exercise 1.15** A subset of $\mathbb{C}^k$ is called *algebraic* if it is the set of solutions of a polynomial equation

$$P(x_1, \ldots, x_k) = 0.$$ 

Show that whenever $A \subseteq \mathbb{C}^k$ is definable by a first order formula in the language of field theory using parameters from $\mathbb{C}$, then $A$ is a Boolean combination of algebraic sets.

What does this say about the expressive power of the language of field theory, in e.g. relation to complex analysis?
Exercise 1.16  Show that \( \mathcal{Q} \), which is \( \mathbb{Q} \) with the standard algebra and ordering, is a model for \( OF \).
Show that \( \mathcal{Q} \) is a prime model for \( OF \).

Exercise 1.17  Carry out the definition of multiplication in the topological completion for the remaining cases.
Verify that the topological completion is indeed an ordered field, and that the original field \( \mathfrak{A} \) may be embedded into its topological completion.

Exercise 1.18  Show that each polynomial function over a real closed field is \( \epsilon - \delta \)-continuous.

Exercise 1.19  Revisit Exercises 1.14 and 1.15.
Reformulate the results proved there to theorems about Real Closed Fields, and prove them.

Exercise 1.20  Prove Theorem 1.4.4 in detail.

Exercise 1.21  Let \( T \) be the complete theory \( DO^+ \) considered in Exercise 1.12.

a) Find a non-principal 1-type as the completion of some infinite set of formulas.

b) **Hard** Describe all the 1-types and show that there is exactly one non-principal 1-type.

c) Which of the three models of \( DO^+ \) is the saturated one, which is the one that is weakly saturated but not saturated and which is the one omitting the one non-principal 1-type?

Exercise 1.22  Let \( DO! \) be the theory introduced on page 36. Show that there is a countable model \( \mathfrak{A} \) for \( DO! \) such that all other models can be embedded into \( \mathfrak{A} \).
Why does it follow that \( DO! \) has only countably many \( n \)-types for every \( n \)?

Exercise 1.23  Let \( L \) be the first order language (with equality) with an enumerable set \( \{ P_i \}_{i \in \mathbb{N}} \) of unary predicate symbols.
Let \( T \) be the theory over \( L \) with non-logical axioms
\[
\exists x (\xi_0(x) \land \cdots \land \xi_{n-1}(x))
\]
where each \( \xi_i \) is either \( P_i \) or \( \neg P_i \).

a) Show that for each \( \xi_0, \ldots, \xi_{n-1} \) as above and each number \( k \) we can formulate and prove in \( T \) that there are at least \( k \) different objects satisfying \( \xi_0 \land \cdots \land \xi_{n-1} \).

b) Let \( L_k \) be \( L \) restricted to the predicate symbols \( P_0, \ldots, P_{k-1} \) and let \( T_k \) be the set of sentences in \( L_k \) that are theorems in \( T \).
Show that \( T_k \) is \( \omega \)-categorical.
c) Show that $T$ is complete.

d) Show that $T$ has uncountably many 1-types.

**Exercise 1.24** Fill in the details in the following sketch of a proof of the omitting type theorem for $n$-types in general:

- We may enumerate all ordered $n$-sequences of Henkin constants in a list $\{(c_{1,k}, \ldots, c_{n,k})\}_{k \in \mathbb{N}}$.
- By induction on $k$, we may find a formula $\phi_k \in X$ such that $T \cup \{ \neg \phi_k(c_{1,k}, \ldots, c_{n,k}) \}$ is consistent with all the Henkin constants.

**Exercise 1.25** Prove Lemma 1.6.6.

**Exercise 1.26** Prove the theorem of constants for $\omega$-logic.

Warning: The proof of the theorem of constants in the first order case use a variable $x$ that occurs nowhere in the proof. In this setting, we do not have proofs, only provability. Moreover, if we had proofs, we might risk that all variables are used in the proof. Discuss how we can overcome this obstacle.

**Exercise 1.27** Prove Lemma 1.6.8.

**Exercise 1.28** Prove Lemma 1.6.9.
Chapter 2

Finitary Model Theory

In Chapter 1 we have touched a little bit on the subject of model theory. To a large extent we have been interested in countable structures; implicitly we considered finite structures to be too simple and uncountable structures to be too advanced. The main reason why we did not go on with more model theory is that we need time and space for computability theory.

In some respects, the finite structures are simple, but in other respects they are the only interesting mathematical structures there are. For instance, a database is normally a finite, but dynamic, structure. Designing logical systems that are adequate for database theory is a challenging task.

In this chapter we will offer the reader a glimpse of a growing field of logic, finite model theory. We are going to prove one of the classical results, the so called 0-1-law. In a sense, this result tells us that first order logic without constants and function symbols have remarkably restricted expressive power, and that we need higher order logic to say something intelligent about a finite structure.

This chapter will have three sections. In the first section we will prove the 0-1-law, in the second section we will extend the language to a 2. order language, and see how the expressive power increases. The final section will be the set of exercises.

2.1 The 0-1-law

For the rest of this chapter, L will be a fixed first order language without constant symbols, without function symbols, and with finitely many predicate symbols.

Definition 2.1.1  a) A complete open description will be an open formula \( M(x_1, \ldots, x_n) \) that is a maximal consistent conjunction of literals in the variables \( x_1, \ldots, x_n \) displayed.
b) If \( M(x_1, \ldots, x_n) \) and \( N(x_1, \ldots, x_n, y_1, \ldots, y_k) \) are complete open descriptions, then \( N \) is an extension of \( M \) if \( \models N \rightarrow M \).

Sometimes it is convenient to use an alternative notation for long conjunctions and disjunctions. In the literature, the notations like \( \wedge_{i \in I} \phi_i \) and \( \vee_{i \in I} \phi_i \) are sometimes used to form infinitary formulas where conjunction and disjunctions are taken over infinitely many subformulas. We will only use the notation when the abbreviations \( \phi_{i_1} \land \cdots \land \phi_{i_k} \) and \( \phi_{i_1} \lor \cdots \lor \phi_{i_k} \) are a bit cumbersome. We will also use the notation \( \vec{x} \) to mean the list of variables \( x_1, \ldots, x_n \), and \( \exists \vec{x} \) will be an abbreviation for \( \exists x_1 \cdots \exists x_n \).

When variables are displayed in a formula like \( M(\vec{x}) \) it is understood that there will be no free variables that are not on display.

With these conventions in place, we will let \( T \) be the theory over \( L \) where the non-logical axioms are all formulas of the form

\[
\forall \vec{x}((\bigwedge_{i \neq j} x_i \neq x_j \land M(\vec{x})) \rightarrow \exists y(\bigwedge_i y \neq x_i \land N(\vec{x}, y)))
\]

whenever \( M \) is an open complete description and \( N \) is an open complete description extending \( M \).

**Theorem 2.1.2** Let \( L \) and \( T \) be as above. Then \( T \) is consistent, has no finite models, and is \( \omega \)-categorical.

**Proof**  
We first prove that \( T \) is consistent, by constructing a model \( \mathfrak{A} \). The idea of the proof is to ad one element at each step, and at each step make one instance of one of the axioms true in the final structure.

Each step is simple, we take an object out of the blue and just make the instance under consideration true by brute force. It is the way we organize the steps in order to ensure that all instances from the final model of all axioms are taken care of that may seem complicated to the reader. Readers with experience from such arguments (similar to, but actually simpler than, our construction of a saturated model for a theory with only countably many element types) may skip parts of this argument.

We will construct an increasing sequence \( \{\mathfrak{A}_k\}_{k \in \mathbb{N}} \) of structures with domain \( \{0, \ldots, k\} \).

Let \( \mathfrak{A}_0 \) be the structure on \( \{0\} \) where \( R^{\mathfrak{A}_0}(0, \ldots, 0) \) holds for all relation symbols \( R \) in \( L \).

Let \( \{(M_i, N_i)\}_{i \in \mathbb{N}} \) be an enumeration of all pairs of complete open descriptions \( M(\vec{x}) \) and \( N(\vec{x}, y) \) such that \( N \) extends \( M \).

Let \( n_i \) be the arity of \( \vec{x} \) in \( M_i \).

For each \( i \), let \( \{\vec{a}_{i,j}\}_{j \in \mathbb{N}} \) be an enumeration of all non-repeating sequences of natural numbers where the sequence has length \( n_i \). We may assume that each component \( a_{m,j,i} \) is bounded by \( j \).

Then at step \( k = 2^i(2j + 1) \) we know that \( \vec{a}_{i,j} \) is a non-repeating sequence from
the domain of $\mathfrak{A}_{k-1}$, so $M_i(\bar{a}_{j,i})$ is a legal instance, with truth value $\bot$ or $\top$. If $M_i(\bar{a}_{j,i})$ is false, we extend $\mathfrak{A}_{k-1}$ to $\mathfrak{A}_k$ by adding $k$ to the domain and making all new instances of atomic statements true.

If $M_i(\bar{a}_{j,i})$ is true, we extend $\mathfrak{A}_{k-1}$ to $\mathfrak{A}_k$ by adding $k$ to the domain and interpreting new instances of atomic formulas such that $N(\bar{a}_{j,i}, k)$ will hold. This will be possible since $N$ is an extension of $M$. If $N$ does not tell us how to interpret an instance of an atomic statement, let it be true.

In this way we will ensure that all the axioms of $T$ will be valid in the limit structure $\mathfrak{A}$ on $\mathbb{N}$, and we do have a model. Thus $T$ is consistent.

Now assume that $\mathfrak{B}$ is a finite model of $T$ with a non-repeating enumeration $\{b_1, \ldots, b_n\}$ of its domain. Let $M(x_1, \ldots, x_n)$ be the unique (up to equivalence) complete open description satisfied by $b_1, \ldots, b_n$.

$M$ will have an extension to a complete open description $N(x_1, \ldots, x_n, y)$.

Since $\mathfrak{B} \models T$ we have that

$$\mathfrak{B} \models \exists y \left( \bigwedge_i y \neq b_i \land M(b_1, \ldots, b_n, y) \right),$$

which is impossible by the choice of $b_1, \ldots, b_n$.

Finally we will show that $T$ is $\omega$-categorical. If $\mathfrak{A}$ and $\mathfrak{B}$ are two countable models for $T$, we will construct an isomorphism using a back and forth construction like the ones we have seen before.

A partial isomorphism will, in this proof, be a map $p : A_0 \to B_0$ where $A_0 \subseteq A$, $B_0 \subseteq B$, $p$ is 1-1 and onto and such that whenever $a_1, \ldots, a_n$ are in $A_0$ and $R$ is an $n$-ary relation symbol, then $R^\mathfrak{A}(a_1, \ldots, a_n)$ if and only if $R^\mathfrak{B}(p(a_1), \ldots, p(a_n))$.

Given $a_1, \ldots, a_n$ there is a unique complete open description $M(x_1, \ldots, x_n)$ (unique up to equivalence) such that $M(a_1, \ldots, a_n)$ holds, and then we may rephrase the definition of a partial isomorphism to:

If $A_0 = \{a_1, \ldots, a_n\}$, $B_0 = \{p(a_1), \ldots, p(a_n)\}$ and $M$ is the complete open description of $\{a_1, \ldots, a_n\}$ then $M$ is also the complete open description of $\{p(a_1), \ldots, p(a_n)\}$. Let $p$ and $a \not\in A_0$ be given. Let $M$ be the complete open description of $A_0$ described above, and let $N$ be the complete open description of $A_0 \cup \{a\}$. The axiom for $M$ and $N$ just implies that $p$ can be extended to a $q$ defined on $a$.

By symmetry we also have that if $p : A_0 \to B_0$ is a partial isomorphism and $b \in B$, then $p$ can be extended to a partial isomorphism $q$ such that $b$ is in the range of $q$.

Then, as before, we can construct an isomorphism between $\mathfrak{A}$ and $\mathfrak{B}$ by piece-wise extensions. This ends the proof.

We will now restrict our attention to $L$-structures $\mathfrak{A}$, where the domain $A$ is a set $\{0, \ldots, n-1\}$. Then the number of $L$-structures will be finite, let $k(L, n)$ denote this number.

Definition 2.1.3 Let $\phi$ be a sentence in $L$. 

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a) Let \( k(\phi, L, n) \) be the number of \( L \)-structures over \( \{0, \ldots, n-1\} \) satisfying \( \phi \) and let the \( n \)-probability of \( \phi \) be
\[
\mu_{L,n}(\phi) = \frac{k(\phi, L, n)}{k(L, n)}.
\]

b) Let the asymptotic probability of \( \phi \) be
\[
\mu_L(\phi) = \lim_{n \to \infty} \mu_{L,n}(\phi)
\]
provided the limit exists.

We will drop the subscript \( L \) from now on.

**Lemma 2.1.4** If \( \phi \) is an axiom of \( T \), then \( \mu(\phi) = 1 \).

**Proof**
Let \( \phi \) be the axiom
\[
\forall \vec{x}(\bigwedge_{i \neq j} x_i \neq x_j \land M(\vec{x}) \rightarrow \exists y(\bigwedge_i y \neq x_i \land N(\vec{x}, y)))
\]
where \( \vec{x} = x_1, \ldots, x_m \).
Assume that \( m < n \). We will estimate \( \mu_n(\neg \phi) \) from above, i.e.
\[
a_n = \mu_n(\exists x_1 \cdots \exists x_m(\bigwedge_{i \neq j} x_i \neq x_j \land M(\vec{x}) \land \forall y(\bigwedge_i y \neq x_i \rightarrow \neg N(\vec{x}, y))))
\]
We will consider formulas where we have used elements from \( \{1, \ldots, n\} \) as instances, without making too much fuss about how to do this strictly according to the book. The local probability of a sentence of this kind is defined in the obvious way.

There are \( n(n-1) \cdots (n-m+1) \) ways of selecting distinct elements for \( x_1, \ldots, x_m \), and by symmetry they are equally probable. Thus we have that
\[
a_n = n(n-1) \cdots (n-m+1) \cdot \mu_n(M(0, \ldots, m-1) \land \forall y(\bigwedge_i y \neq i \rightarrow \neg N(0, \ldots, m-1, y)))
\]
which is bounded by
\[
b_n = n^m \cdot \mu_n(M(0, \ldots, m-1) \land \forall y(\bigwedge_i y \neq i \rightarrow \neg N(0, \ldots, m-1, y)))
\]
which again is bounded by
\[
c_n = n^m \cdot \mu_n(\forall y(\bigwedge_{i=0}^{m-1} y \neq i \rightarrow \neg N(0, \ldots, m-1, y))).
\]
Let $\xi_1, \ldots, \xi_l$ be all the literals in $N$ with an occurrence of the variable $y$. Then

$$c_n = n^m \cdot \mu_n(\forall y(\bigwedge_{i=0}^{m-1} y \neq i \rightarrow \bigvee_{j=1}^l \neg \xi_j(0, \ldots, m - 1, y))).$$

For each $k$ with $m \leq k \leq n$ and $j \leq l$, the probability of $\xi_j(0, \ldots, m - 1, k)$ is exactly $\frac{1}{2}$. Thus we may evaluate

$$c_n = n^m \cdot \prod_{i=m}^n \mu_n(\bigvee_{j=1}^l \neg \xi_j(0, \ldots, m - 1, i)) = n^m(1 - \frac{1}{2^l})^{n-m}.$$

Since $m$ and $l$ are fixed and $(1 - \frac{1}{2^l}) < 1$, we have that

$$\lim_{n \to \infty} n^m(1 - \frac{1}{2^l})^{n-m} = 0.$$

This proves the theorem.

**Corollary 2.1.5** Let $L$ and $T$ be as above, $\phi$ a sentence in $L$. Then $\mu(\phi) = 1$ or $\mu(\phi) = 0$.

**Proof**

Since $T$ is $\omega$-categorical, we know in particular that $T$ is complete. If $T \vdash \phi$, there are axioms $\phi_1, \ldots, \phi_k$ in $T$ such that

$$\phi_1, \ldots, \phi_k \vdash \phi.$$ 

Since

$$\mu_n(\phi_1 \land \cdots \land \phi_k) \geq 1 - \sum_{i=1}^k (1 - \mu_n(\phi_i))\,$$

it follows that $\mu(\phi) = 1$. By the same argument, if $T \vdash \neg \phi$ we have that $\mu(\phi) = 0$ since $\mu(\neg \phi) = 1$.

**An application to graph theory**

In this section we will let a **directed graph** simply be a binary relation $R$ on a set $G$. In order to get the feeling of a graph-theoretical environment, we will call the ordered pairs $(a, b)$ for **oriented edges** with tail $a$ and head $b$, and we call the elements of $G$ for **nodes**. Graph theory also covers the case of non-oriented edges, the possibility of several edges between two nodes, no edges from a node to itself etc. For all these definitions of graphs, a similar phenomenon as the one we will discuss takes place.

A **path** from $a$ to $b$ will be a finite sequence of edges $e_1, \ldots, e_n$ such that

- $a$ is the head of $e_1$.
- $b$ is the tail of $e_n$. 

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For all \( i < n \), the head of \( e_i \) is the tail of \( e_{i+1} \).

A path like this will have length \( n \). We say that the graph is \textit{connected} if there is a path from \( a \) to \( b \) for all nodes \( a \) and \( b \) (this concept is more commonly used for undirected graphs, so do not use it in a context independent of this section). We will call a graph \textit{intimately connected} (not a standard concept in any context) if all nodes \( a \) and \( b \) are connected via a path of length at most 2.

**Corollary 2.1.6** Let \( P(n) \) be the probability that an arbitrary graph on \( \{0, \ldots, n\} \) is intimately connected. Then

\[
\lim_{n \to \infty} P(n) = 1.
\]

**Proof**

Let \( L \) be the language of graph theory, and let \( T \) be as above. Consider the sentence

\[
\phi = \forall x \forall y \exists z (R(x, z) \land R(z, y)).
\]

\( \phi \) will be a theorem in \( T \), and thus have asymptotic probability 1.

This example shows that most finite graphs are non-interesting. In an effort to make this corollary sound interesting, we claim that it demonstrates that one must show some care when describing the virtue of probabilistic algorithms. The probabilistic algorithm ‘Yes’ solves the problem ‘is \( G \) connected?’ with better and better probability the larger the finite graph \( G \) is, but it is of course absolutely worthless.

In Exercise 2.2 we show that most transitive finite binary relations are rather non-interesting.

### 2.2 Second order languages

We have shown that the first order theory of algebraically closed fields of a fixed characteristic is complete. We have shown that the natural numbers cannot be characterized up to isomorphism by any first order theory. In this chapter we have shown the 0-1-law of finite model theory for first order formulas. These and other results indicate that first order languages are poor in expressive power.

These languages are called \textit{first order} because we only accept quantifiers over the domain of the structure. A standard mathematical format of the induction axiom for number theory is as follows:

\[
\forall A(0 \in A \land \forall x \in \mathbb{N}(x \in A \rightarrow x + 1 \in A) \rightarrow \forall x \in \mathbb{N}(x \in A)).
\]

Here we let \( A \) vary over all sets, in particular over all subsets of \( \mathbb{N} \). Quantifiers ranging over all subsets of a domain, or more generally, over all predicates of some fixed arity on a domain, will be called \textit{second order}. We will be more precise in a while, but notice already now that the natural numbers can be described up to isomorphism by second order formulas, see Exercise 2.4.
One problem with interpreting a second order quantifier is that the interpretation is not absolute. By this we mean that the powerset of an infinite set is not fully understood. There are second order sentences in number theory where the truth value will depend on chosen axiomatisations of set theory. This will be discussed at more depth in a course on axiomatic set theory. If we restrict our attention to finite structures, the same argument against second order logic in general is not valid. Given a finite set we have, at least theoretically, a full control of the power set.

We will show that the expressive power of second order quantifiers is such that the 0-1-law does not hold any more. A more systematical treatment of second order logic and other extensions of first order logic is beyond the scope of this text.

**Definition 2.2.1**

a) A second order language $L^2$ will be a first order language $L$ extended with variables $X^k_n$ for predicates of arity $k$.

b) If $t_1,\ldots,t_k$ are terms in $L$, then $X^k_n(t_1,\ldots,t_k)$ is a new atomic formula.

c) The class of second order formulas is closed under boolean connectives and first order quantifiers in the same way as the first order formulas are.

d) If $\phi$ is a second order formula, then $\exists X^k_n \phi$ and $\forall X^k_n \phi$ are second order formulas.

We will not go through the full definition of how second order formulas are interpreted over an $L$-structure. The point is that they can be interpreted, and that if $\mathfrak{A}$ and $\mathfrak{B}$ are isomorphic $L$-structures, then they will satisfy the same second order sentences.

From now on we will restrict ourselves to the first order language $L_\text{eq}$ of equality, and its second order extension $L^2_\text{eq}$. Then there is exactly one structure with domain $\{0,\ldots,n-1\}$ for each $n$.

**Theorem 2.2.2** *There is a second order sentence that is true exactly for the finite structures with an even number of elements.*

**Proof**

In the sentence below, $X$ will be a predicate variable of arity 1 and $Y$ will be a predicate variable of arity 2.

Let us first look at the sentence, and discuss why it works afterwards:

$$\exists X \exists Y [\forall x \exists y Y(x,y) \land \forall y \exists x Y(x,y)$$

$$\land \forall x \forall y \forall z (Y(x,y) \land Y(x,z) \rightarrow y = z)$$

$$\land \forall x \forall y \forall z (Y(x,y) \land Y(z,y) \rightarrow x = z)$$

$$\land \forall x \forall y (Y(x,y) \rightarrow (X(x) \leftrightarrow \neg X(y))).$$

The first and third conjuncts express that $Y$ is the graph of a function $f_Y$. The second and fourth conjuncts express that $f_Y$ is onto and one-to-one, i.e.
a bijection. The last conjunct expresses that \( f_Y \) is a bijection between \( X \) and the complement of \( X \). If the domain of the structure is finite, we may find interpretations of \( X \) and \( Y \) satisfying this exactly when the total number of elements is an even one.

Clearly the sentence above does not satisfy the 0-1-law, so we have proved

**Corollary 2.2.3** The 0-1-law is not satisfied in general by second order sentences.

In the proof of Theorem 2.2.2 we saw that we can express that \( Y \) is the graph of a function and other properties of functions using first order expressions in the variable \( Y \). Other properties that we may express is

1. \( X^1 \) is finite.
2. \( X^1 \) is infinite.
3. \( \langle X^1, Y^2 \rangle \) is a well ordering
4. A ring is Noetherian (which means that any decreasing sequence of ideals is finite)

The verifications of the first three facts are left as Exercise 2.5. The proof of Fact 4. is easy if you know the definitions of Noetherian rings and ideals, but will be very hard if you do not know these definitions.

### 2.3 Exercises to Chapter 2

**Exercise 2.1** Show that there is an upper bound on the length of a complete open description \( M(x_1, \ldots, x_n) \) depending on \( n \) and the signature of the language, and that whenever \( M(x_1, \ldots, x_n) \) is a complete open description and \( y_1, \ldots, y_k \) are new variables, then \( M \) has an extension to a complete open description \( N(x_1, \ldots, x_n, y_1, \ldots, y_m) \).

**Exercise 2.2** In automata theory, in logic and also in other branches of mathematics or theoretical computer science we sometimes define a relation by considering the reflexive and transfinite closure of a one-step-relation (e.g. one-step action of a Turing machine or one-step deduction in a formal theory). The transitive closure of a binary relation cannot be axiomatized in first order logic (if you did not know this, consider it as a part of this exercise to prove it) and is defined by

\[
R^*(a, b) \iff \exists n \in \mathbb{N} \exists a_1, \ldots, a_n (a = a_1 \land b = a_n \land \forall i < n R(a_i, a_{i+1})).
\]

Let \( R \) be an arbitrary binary relation on \( \{0, \ldots, n\} \). Show that with asymptotic probability 1, \( R^*(i, j) \) for all \( i, j \leq n \).

What are the asymptotic probabilities of the interesting properties
R is reflexive?
R is nowhere reflexive?
R is transitive?
R is symmetric?

Exercise 2.3 Let \( L \) be the language with \( = \) as the only symbol.

a) Show that for every sentence \( \phi \) in \( L \) and every \( n \) we have that \( \mu_n(\phi) \in \{0, 1\} \).

b) Show that the function \( n \mapsto \mu_n(\phi) \) will be eventually constant for each sentence \( \phi \).

c) Evaluate from where the function \( n \mapsto \mu_n(\phi) \) is eventually constant, in terms of some number related to \( \phi \), and use this to show that the pure theory of equality is decidable.

Exercise 2.4 Find a set of first and second order sentences in the language of number theory that will have isomorphic copies of \( \mathbb{N} \) as the only models.

Exercise 2.5 A standard set-theoretical characterization of infinity is that a set \( A \) is infinite exactly when there is a bijection between \( A \) and a proper subset of \( A \). Use this to show that we may express that a set is finite by a second order formula.

Find a second order formula that express that \( Y^2 \) is a well ordering of \( X^1 \).

If you want to go on with harder problems, we suggest that you continue with Exercise 5.4.
Chapter 3

Classical Computability Theory

3.1 The foundation, Turing’s analysis

In Leary [1] the recursive functions are defined as those that can be represented in elementary number theory. \( f : \mathbb{N}^k \to \mathbb{N} \) is recursive if there is a formula \( \phi(x_1, \ldots, x_k, y) \) such that for all \( n_1, \ldots, n_k, m \) we have that \( f(n_1, \ldots, n_k) = m \) if and only if

\[ N \vdash \phi(c_{n_1}, \ldots, c_{n_k}, y) \iff y = c_m. \]

Here \( c_n \) is the numeral for \( n \), and \( N \) is elementary number theory.

The advantage of this definition is that it is well suited for proving Gödel’s incompleteness theorem without introducing too many new concepts. The problem is that there is absolutely no conceptual analysis of the notion of computability behind this definition.

Gödel defines a class of recursive functions by recursion (pun intended). His aim is to define a sufficiently rich class for handling algorithms for e.g. substitution of a term for a variable, and for coding the meta-theory of a formal theory, but sufficiently simple to enable us to show that any recursive function will be definable, and actually, representable as described above. Gödel’s computable functions are now known as the \( \mu \)-recursive ones.

We are going to base our study of computability on an approach due to Kleene, and we are going to restrict ourselves to computable functions defined on the natural numbers. In many respects, computing can be described as manipulation of symbols following a given set of rules. The symbols are not then natural numbers, and different ways of representing natural numbers (binary, decadic, via numerals, Roman figures etc.) might give different concepts of computing with numbers.

The best mathematical model for computability and computations is due to Alan Turing. He defined what is now known as Turing machines, small finite state machines operating on an infinite one-dimensional tape and doing
symbol manipulations on this tape. The machine will, between each step of the computation, be in one of finitely many states. It will read one of the symbols on the tape, and dependent of the state it is in and the symbol on the tape that it reads, it will according to a fixed rule change its state, rewrite the symbol and move to the symbol to the right or to the left. It may of course stay in the same state, it may of course keep the symbol on the tape as it is, and sometimes it may be convenient to let it stay where it is. We think of the tape as consisting of squares, like an old fashioned movie-tape.

A Turing machine $M$ is determined by

1. A finite alphabet $\Sigma$ including one special symbol $Bl$ for an empty square of the tape.
2. A finite set $K$ of states, with one special state $s \in K$ called the initial state.
3. A partial function $\delta : (\Sigma \times K) \rightarrow (\Sigma \times K \times \{L, S, R\})$, where $L$ means "left", $R$ means "right" and $S$ means "stay where you are".

In the literature you will find many variations in the definition of a Turing machine, but they all have the same computational power. In this exposition, we decided to let the function $\delta$, which rules the action of the machine, be partial. If $\delta(\sigma, p)$ is undefined, the machine will be in a halting situation, which means that the computation comes to an end. We are not going to give a precise definition of the operational semantics for Turing machines, but are content with an intuitive description:

The operational semantics of Turing Machines

Let $M = (\Sigma, K, s, \delta)$ be a Turing machine. The starting configuration of $M$ will consist of a word $w$ (called the input) in $\Sigma$ written on an otherwise blank tape that is infinite in both directions, a position on the tape and the initial state $s$.

At each step, the machine $M$ may enter a new state, rewrite the symbol at its position on the tape and shift its position one square to the right or to the left, all according to the function $\delta$.

If $M$ is in stage $p$ and reads the symbol $\sigma$, and $\delta(\sigma, p)$ is undefined, $M$ halts. Then we who observe $M$ will know this, and we will be able to read the content of the tape, which will be called the output.

Normally there will be some conventions for how to organize the input configuration, e.g. that there should be no blanks in the input word and that the machine will be started at the first blank square to the left or to the right of the input word. Then the following makes sense:

**Definition 3.1.1** Let $\Sigma_0$ be an alphabet not containing the symbol $Bl$

Let $\Sigma_0^*$ be the set of finite words over $\Sigma_0$.

Let $f : \Sigma_0^* \rightarrow \Sigma_0^*$ be a partial function.

We say that $f$ is Turing computable if there is a Turing machine $M$ over an
alphabet $\Sigma \supset \Sigma_0$ such that if $M$ is started with $w \in \Sigma^*$ on the tape, then it halts if and only if $f(w)$ is defined, and then with $f(w)$ as the output word.

Turing claimed that a Turing machine can

- Search systematically for pieces of information.
- Remember pieces of information
- Rewrite contents

all according to fixed rules. As an example we may consider the map $n \mapsto n!$ and the map $n \mapsto m^m$. We all agree that these maps are in principle computable, but try to think about how a computation of $10^{101}$! could be carried out: We would need a lot of book-keeping devices in order to be at the top of the situation at each stage, but nothing that is not covered by the three items above.

We follow Turing in claiming that his model is a good mathematical model for algorithmic computations. We are, however, not going to make this compendium into a study of Turing machines. The interested reader should consult other textbooks on the subject.

The two key results are:

**Theorem 3.1.2** There is a fixed alphabet $\Sigma$ such that for any alphabet $\Sigma'$, any Turing machine $M$ over $\Sigma'$ may be coded as a word $[M]$ in $\sigma$ and every word $w$ in $\Sigma'$ may be coded as a word $[w]$ in $\Sigma$ such that the partial function

$$U([M][w]) = [M(w)]$$

is Turing computable, where we write $M(w)$ for the output word if the input word is $w$.

This is known as the existence of a *Universal Turing Machine*.

The *Halting Problem* is the following:

Given a Turing machine $M$ and an input $w$, will $M$ eventually come to a halt when started on $w$?

**Theorem 3.1.3** There is no Turing machine $H$ that solves the Halting Problem in the following sense:

$H([M][w])$ will always halt, and will halt with an empty output word if and only if $M$ halts on $w$.

Our approach to computability will be more in the original style of Gödel, we will study functions defined for natural numbers only. However, the results that we obtain will be relevant for the more general approaches as well. This is based on what is known as the Church-Turing Thesis, which we phrase like this:

*All algorithms can be simulated by a Turing Machine,*

and the fact that Turing-computability can be reduced to the notion we will be working with. This is left as one of the minor projects in Chapter 5, see Exercise 5.6.
3.2 Computable functions and c.e. sets

3.2.1 The primitive recursive functions

The basic definition

Recursion means ‘backtracking’, and in pre-Church/Kleene mathematics the term *recursive function* was used for the functions defined by iterated recursion. In this pre-Church/Kleene context a recursive definition will be a definition of a function on the natural numbers, where we give one initial value \( f(0) \), and define \( f(k+1) \) as a function of \( f(k) \) and \( k \). E.g. Skolem used the term ‘recursive function’ in this way. Following Kleene, we will call these functions *primitive recursive*.

We let \( \vec{x} \) and \( \vec{y} \) etc. denote ordered sequences of natural numbers of some fixed length. Normally the length will not be specified, but will be clear from the context.

**Definition 3.2.1** *The primitive recursive functions* \( f : \mathbb{N}^n \to \mathbb{N} \) will be the least class of functions satisfying:

i) \( f(x, \vec{y}) = x + 1 \) is primitive recursive.

ii) \( I_{1,n}(x_1, \ldots, x_n) = x_i \) is primitive recursive.

iii) \( f(\vec{x}) = q \) is primitive recursive for each \( q \in \mathbb{N} \).

iv) If \( g \) is \( n \)-ary and primitive recursive, and \( f_1, \ldots, f_n \) are \( m \)-ary and primitive recursive, then the composition

\[
h(\vec{x}) = g(f_1(\vec{x}), \ldots, f_n(\vec{x}))
\]

is primitive recursive.

v) If \( g \) and \( h \) are primitive recursive of arity \( n \) and \( n + 2 \) resp., then \( f \) is primitive recursive where

\[
\begin{align*}
f(0, \vec{y}) &= g(\vec{y}) \\
f(x + 1, \vec{y}) &= h(f(x, \vec{y}), x, \vec{y})
\end{align*}
\]

We say that \( f \) is defined by *primitive recursion* or by the recursion operator from \( g \) and \( h \).

The pool of primitive recursive functions and sets

Standard number theoretical functions like addition, multiplication, exponentiation and factorial will all be primitive recursive. For instance \( f(x, y) = x + y \) may be defined by

- \( x + 0 = x \)
- \( x + (y + 1) = (x + y) + 1 \).
Subtraction is not primitive recursive for the simple reason that it leads us outside \( \mathbb{N} \). We define a modified subtraction. This will be primitive recursive, see Exercise 3.1.

**Definition 3.2.2** We let \( \cdot \) be the modified subtraction defined by
\[
\begin{align*}
\cdot - y &= \cdot x - y \text{ if } y \leq x \\
\cdot - y &= 0 \text{ if } x \leq y.
\end{align*}
\]

By the definition of the primitive recursive functions, all projection maps
\[
I_{i,n}(x_1, \ldots, x_n) = x_i
\]
are primitive recursive. We can use this to consider any function of a set of variables as a function of a larger set of variables, as in
\[
f(x_1, x_2, x_3) = g(I_{1,3}(x_1, x_2, x_3), I_{3,3}(x_1, x_2, x_3)) = g(x_1, x_3).
\]

Thus we will not need to be concerned with the requirement that every function in a composition must be of the same arity. This will simplify some of the descriptions of primitive recursive functions.

**Definition 3.2.3** Let \( A \subseteq \mathbb{N}^n \). The characteristic function of \( A \) will be the function
\[
K_A : \mathbb{N}^n \to \mathbb{N}
\]
that is 1 on \( A \) and 0 outside \( A \).

**Definition 3.2.4** A set \( A \subseteq \mathbb{N}^n \) is primitive recursive if the characteristic function \( K_A \) is primitive recursive.

The primitive recursive sets will form a Boolean algebra for every dimension, see Exercise 3.2.

If \( A \subseteq \mathbb{N}^k \) is primitive recursive and \( f, g : \mathbb{N}^k \to \mathbb{N} \) are primitive recursive, we may define a primitive recursive function \( h \) by cases as follows:
\[
h(\bar{x}) = K_A(\bar{x}) \cdot f(\bar{x}) + (1 - K_A(\bar{x})) \cdot g(\bar{x}).
\]

We see that \( h(\bar{x}) = f(\bar{x}) \) when \( \bar{x} \in A \) and \( h(\bar{x}) = g(\bar{x}) \) otherwise.

Every set defined from = and < using propositional calculus will be primitive recursive. We may also use functions known to be primitive recursive in the definition of primitive recursive sets, and primitive recursive sets when we describe primitive recursive functions. We will leave the verification of most of this to the reader, just state the properties of primitive recursion that are useful to us. The proofs are simple, and there is no harm leaving them as exercises. We will however prove one basic (and simple) lemma:
Lemma 3.2.5 Let $f : \mathbb{N}^{1+n} \rightarrow \mathbb{N}$ be primitive recursive. Then the function $g$ defined as the bounded product

$$g(x, \vec{y}) = \prod_{z \leq x} f(z, \vec{y})$$

will be primitive recursive.

Proof

We define $g$ by primitive recursion as follows

$$g(0, \vec{y}) = f(0, \vec{y})$$

$$g(x + 1, \vec{y}) = g(x, \vec{y}) \cdot f(x + 1, \vec{y})$$

By the same argument we can show that the primitive recursive functions will be closed under bounded sums. What is more important to us is that the primitive recursive sets will be closed under bounded quantification. This is an important strengthening of the language we may use to define primitive recursive sets and functions.

Lemma 3.2.6 Let $A \subseteq \mathbb{N}^{1+n}$ be primitive recursive. Then the following sets are primitive recursive

a) $B = \{(x, \vec{y}) ; \exists z \leq y((z, \vec{y}) \in A)\}$

b) $C = \{(x, \vec{y}) ; \forall z \leq y((z, \vec{y}) \in A)\}$

The proof is left as Exercise 3.3.

In computability theory, the $\mu$-operator is important. Within primitive recursion theory we may often use bounded search, or a bounded $\mu$-operator:

Lemma 3.2.7 Let $f : \mathbb{N}^{1+n} \rightarrow \mathbb{N}$ be primitive recursive. Then

$$g(x, \vec{y}) = \mu_{<x} z. (f(z, \vec{y}) = 0)$$

is primitive recursive, where $g(x, \vec{y})$ is the least $z$ such that $f(z, \vec{y}) = 0$ if there is one such $z < x$, while $g(x, \vec{y}) = x$ otherwise.

$g$ can be defined using primitive recursion and definition by cases. The details are left as Exercise 3.4.

The interplay between the primitive recursive functions and the primitive recursive sets is ruled by the following principles:

Theorem 3.2.8 a) Every set definable from $=$ and $<$ using primitive recursive functions, boolean operators and bounded quantifiers will be primitive recursive.

b) Every function defined by the schemes of primitive recursion, bounded search over a primitive recursive set and definition by cases over a finite partition of $\mathbb{N}^n$ into primitive recursive sets will be primitive recursive.

There is no need to prove this theorem, since it in a sense is the synthesis of what has been proved so far. The consequence is the level of freedom in defining primitive recursive functions and relations we have obtained. This freedom will be sufficient when we later claim that certain facts are trivial to prove.
Sequence numbers

One important use of primitive recursion is the coding of finite sequences. Gödel needed an elaborate way of coding finite sequences via the so called $\beta$-function. As we mentioned above, it was important for Gödel to show that the recursive functions are definable in ordinary number theory. Since this is not important to us to the same extent, we will use full primitive recursion in coding such sequences. All proofs of the lemmas below are trivial and can safely be left for the reader.

Lemma 3.2.9

a) The set of prime numbers (starting with 2 as the least prime number) is primitive recursive.

b) The monotone enumeration $\{p_i\}_{i \in \mathbb{N}}$ of the prime numbers is primitive recursive (with $p_0 = 2$).

We now define the sequence numbers:

Definition 3.2.10

Let $x_0, \ldots, x_{n-1}$ be a finite sequence of numbers, $n = 0$ corresponding to the empty sequence.

We let the corresponding sequence number $\langle x_0, \ldots, x_{n-1} \rangle$ be the number

$$2^{\text{lh}(y)} \prod_{i=0}^{n-1} p_{x_i+1}^{x_i}$$

If $y = \langle x_0, \ldots, x_{n-1} \rangle$, we let $lh(y)$ (the length of $y$) be $n$ and $(y)_i = x_i$.

If $x = \langle x_0, \ldots, x_{n-1} \rangle$ and $y = \langle y_0, \ldots, y_{m-1} \rangle$, we let $x * y$ be the sequence number of the concatenation, i.e.

$$x * y = \langle x_0, \ldots, x_{n-1}, y_0, \ldots, y_{m-1} \rangle.$$

Sometimes, when it is clear from the context, we will misuse this terminology and let

$$a * \langle x_0, \ldots, x_{n-1} \rangle = \langle a, x_0, \ldots, x_{n-1} \rangle.$$

Lemma 3.2.11

a) The set of sequence numbers is primitive recursive, and the sequence numbering is one-to-one (but not surjective).

b) The function $lh$ is the restriction of a primitive recursive function to the set of sequence numbers.

c) The function $\text{coor}(y, i) = (y)_i$ (the $i$'th coordinate of $y$) defined for sequence numbers $y$ of length $> i$ is the restriction of a primitive recursive function of two variables.

d) For each $n$, the function

$$\text{seq}_n(x_0, \ldots, x_{n-1}) = \langle x_0, \ldots, x_{n-1} \rangle$$

is primitive recursive.
e) For all sequences \(x_0, \ldots, x_{n-1}\) and \(i < n\), \(x_i < \langle x_0, \ldots, x_{n-1}\rangle\)

f) \(1\) is the sequence number of the empty sequence.

h) Concatenation of sequence numbers is primitive recursive.

It makes no sense to say that the full sequence numbering is primitive recursive. However any numbering satisfying Lemma 3.2.11 can be used for the purposes of this course, so there is no need to learn the details of this definition. Occasionally it may be useful to assume that the sequence numbering is surjective. This can be achieved, see Exercise 3.5.

Primitive recursion is technically defined on successor numbers via the value on the predecessor. Sometimes it is useful to use the value on all numbers less than the argument \(x\) in order to define the value on \(x\). In order to see that this is harmless, we use the following construction and the result to be proved in Exercise 3.6:

**Definition 3.2.12** Let \(f : \mathbb{N} \to \mathbb{N}\).
Let \(\bar{f}(n) = \langle f(0), \ldots, f(n-1) \rangle\).

We will sometimes use an alternative coding of pairs that is both 1-1 and onto:

**Definition 3.2.13** Let \(P : \mathbb{N}^2 \to \mathbb{N}\).

Let \(\pi_1\) and \(\pi_2\) be the two projection maps such that for any \(x\)

\[P(\pi_1(x), \pi_2(x)) = x.\]

\(P, \pi_1\) and \(\pi_2\) are primitive recursive. The verifications are left for the reader as Exercise 3.5 e).

Ackermann proved that there is a total computable function that is not primitive recursive. His observation was that in the list of functions

\[f_0(x, y) = x + 1, \quad f_1(x, y) = x + y, \quad f_2(x, y) = xy, \quad f_3(x, y) = x^y\]

each function but the first is defined as a \(y\)-iteration of the previous one.

For \(n \geq 3\) we may then define

\[f_n(x, 0) = 1, \quad f_n(x, y + 1) = f_{n-1}(f_n(x, y), x)\]

This defines the generalized exponentiations, or the Ackermann-branches. The three-place function

\[f(n, x, y) = f_n(x, y)\]

is not primitive recursive.

We will not prove Ackermann’s result as stated, but in Exercise 3.7 we will see how we may define a function using a double recursion that is not primitive recursive. It will be easier to solve this exercise after the introduction to Kleene indexing.
3.2.2 The computable functions

The $\mu$-operator

We extend the definition of the primitive recursive functions to a definition of the computable functions by adding one principle of infinite search. We will consider the construction

$$g(\vec{x}) = \mu x. f(x, \vec{x}) = 0.$$  

In order to understand this definition, we must discuss what we actually mean, i.e. which algorithm this is supposed to represent.

Intuitively we want $\mu x. g(x, \vec{x}) = 0$ to be the least $x$ such that $g(x, \vec{x})$ is 0. If $g(x, \vec{x}) \in \mathbb{N}$ for all $x$, this is unproblematic, we search for this least $x$ by computing $g(0, \vec{x})$, $g(1, \vec{x})$ and so on until we find one value of $x$ giving 0 as the outcome. Now, if there is no such $x$ we will search in vain, or in more technical terms, our procedure will be non-terminating. This forces us to introduce partial functions, i.e. functions being undefined on certain arguments. This further forces us to be careful about our interpretation of the $\mu$-operator, we may in the search for the least $x$ such that $g(x, \vec{x}) = 0$ face a $z$ for which $g(z, \vec{x})$ is undefined before we find a $z$ for which $g(z, \vec{x}) = 0$. In this case we will let $\mu x. g(x, \vec{x}) = 0$ be undefined, realizing that the algorithm for computing this number that we have in mind will be non-terminating.

With this clarification we ad the following.

**Definition 3.2.14** The **computable functions** is the least class of partial functions $f : \mathbb{N}^n \rightarrow \mathbb{N}$ satisfying i) - v) in the definition of the primitive recursive functions, together with the extra clause

vi) If $g : \mathbb{N}^{n+1} \rightarrow \mathbb{N}$ is computable then

$$f(\vec{x}) = \mu x. g(x, \vec{x}) = 0$$

is computable.

A set is **computable** if the characteristic function is computable.

A **total computable function** is a computable function terminating on all inputs from the domain $\mathbb{N}^n$.

**Remark 3.2.15** Gödel’s $\mu$-recursive functions will be the total computable functions.

We must have in mind that the characteristic function of a set is total, so dealing with computable sets and with total computable functions will be much of the same. This is made precise in the following lemma:

**Lemma 3.2.16** Let $f : \mathbb{N}^n \rightarrow \mathbb{N}$ be total. Then the following are equivalent:

i) $f$ is computable.
ii) The graph of $f$, seen as a subset of $\mathbb{N}^{n+1}$, is computable.

Proof
Let $A$ be the graph of $f$, $K_A$ the characteristic function of $A$.
If $f$ is computable, then

$$K_A(\vec{x}, y) = K_=(f(\vec{x}), y)$$

and $K_=$ is primitive recursive, see Exercise 3.1, g) and e).
If $K_A$ is computable, then

$$f(\vec{x}) = \mu y. 1 - K_A(\vec{x}, y) = 0,$$

so $f$ is computable.

Remark 3.2.17 This result will not hold for primitive recursion, there will be functions with primitive recursive graphs that are not primitive recursive.

There is an interplay between the total computable functions and the computable sets resembling Theorem 3.2.8 as follows

Theorem 3.2.18
a) Any set definable from computable sets and total computable functions using boolean valued operators and bounded quantifiers will be computable.

b) If $A \subseteq \mathbb{N}^{n+1}$ is computable, then $g$ is computable, where $g(\vec{x})$ is the least number $z$ such that $(z, \vec{x}) \in A$.

The proof is trivial.

Kleene’s $T$-predicate
The definition of the computable functions is technically an inductive definition of a class of partial functions. It is, however, important to be aware of the computational interpretation of this definition, the so to say ‘operational semantics’. Every partial computable function is given by a term in a language with constants for each number, the $+1$ function and symbols for the recursion operator and the $\mu$-operator.

This operational semantics then tells us that there are actual computations going on. Now we will define the concept of a computation tree. A computation tree will be a number coding every step in a computation up to the final outcome. In order to be able to do so, we need a Gödel numbering, or an indexing, of the computable functions. Observe that the numbering will not be 1-1, we will enumerate the algorithms or the terms, and then only indirectly enumerate the computable functions.

Definition 3.2.19 For each number $e$ and simultaneously for all $n$ we define the partial function $\phi^n_e$ of $n$ variables $\vec{x} = (x_1, \ldots, x_n)$ as follows:
i) If \( e = \langle 1 \rangle \), let \( \phi^n_e(\vec{x}) = x_1 + 1 \).

ii) If \( e = \langle 2, i \rangle \) and \( 1 \leq i \leq n \), let \( \phi^n_e(\vec{x}) = x_i \).

iii) If \( e = \langle 3 \rangle \), let \( \phi^n_e(\vec{x}) = q \).

iv) If \( e = \langle 4, e', d_1, \ldots, d_m \rangle \) let

\[
\phi_e(\vec{x}) = \phi^n_{e'}(\phi^{n}_{d_1}(\vec{x}), \ldots, \phi^{n}_{d_m}(\vec{x})).
\]

v) If \( e = \langle 5, d_1, d_2 \rangle \) then

\[
\begin{align*}
\phi^{n+1}_e(0, \vec{x}) &= \phi^n_{d_1}(\vec{x}) \\
\phi^{n+1}_e(x + 1, \vec{x}) &= \phi^{n+2}_{d_2}(\phi^{n+1}_e(x, \vec{x}), x, \vec{x}).
\end{align*}
\]

vi) If \( e = \langle 6, d \rangle \) then

\[
\phi^n_e(\vec{x}) = \mu z. \phi^{n+1}_e(z, \vec{x}) = 0
\]

Otherwise if neither of i) - vi) above applies, let \( \phi^n_e(\vec{x}) \) be undefined.

**Remark 3.2.20** We have defined \( \phi_e(\vec{x}) \) for every index \( e \) and every input \( \vec{x} \), either as an undefined value or as a natural number. Indeed, if we get a natural number, this number will be unique, see Exercise 3.8. When no misunderstanding should occur, we will drop the superscript \( n \) and write \( \phi_e(\vec{x}) \) instead of \( \phi^n_e(\vec{x}) \).

**Definition 3.2.21** We write \( \phi_e(\vec{x}) \downarrow \) if there is a \( y \) with \( \phi_e(\vec{x}) = y \). We then say that \( \phi_e(\vec{x}) \) terminates.

We are now ready to use the sequence numbering and this indexing to define computation trees. Each terminating computation will have a unique computation tree, a number coding each step of the computation from the input to the output. We will actually be overloading this code with information, but for our purposes this is harmless. What is important to us is that information retrieval will be easy.

**Definition 3.2.22** Let \( \phi_e(\vec{x}) = y \). By primitive recursion on \( e \) we define the computation tree of \( \phi_e(\vec{x}) = y \) as follows, assuming that the index \( e \) will be the index of the corresponding case:

i) \( \langle e, \vec{x}, x_1 + 1 \rangle \) is the computation tree for \( \phi_e(x, \vec{y}) = x + 1 \).

ii) \( \langle e, \vec{x}, x_i \rangle \) is the computation tree for \( \phi_e(\vec{x}) = x_i \).

iii) \( \langle e, \vec{x}, q \rangle \) is the computation tree for \( \phi_e(\vec{x}) = q \).

iv) \( \langle e, t, t_1, \ldots, t_n, y \rangle \) is the computation tree in this case, where each \( t_i \) is the computation tree for \( \phi_{d_i}(\vec{x}) = z_i \) and \( t \) is the computation tree for \( \phi_{e'}(\vec{x}) = y \).
v) \( \langle e, 0, t, y \rangle \) is the computation tree for \( \phi_e(0, \vec{x}) = y \) when \( t \) is the computation tree for \( \phi_d(\vec{x}) = y \).

\( \langle e, x + 1, t_1, t_2, y \rangle \) is the computation tree for \( \phi_e(x + 1, \vec{x}) = y \) when \( t_1 \) is the computation tree for \( \phi_e(x, \vec{x}) = z \) and \( t_2 \) is the computation tree for \( \phi_d(z, x, \vec{x}) = y \).

vi) \( \langle e, t_0, \ldots, t_{y-1}, t_y, y \rangle \) is the computation tree in this case, where \( t_i \) is the computation tree for \( \phi_d(i, \vec{x}) = z_i \neq 0 \) for \( i < y \) and \( t_y \) is the computation tree for \( \phi_d(y, \vec{x}) = 0 \).

We say that \( t \) is a computation tree for \( \phi^n_e(\vec{x}) \) if for some \( y \), \( t \) is the computation tree for \( \phi^n_e(\vec{x}) = y \).

We are now ready to define Kleene’s \( T \)-predicate:

**Definition 3.2.23** Let

\[
T_n(e, x_1, \ldots, x_n, t)
\]

if \( t \) is a computation tree for \( \phi_e(x_1, \ldots, x_n) \)

We will normally write \( T \) instead of \( T_1 \).

**Theorem 3.2.24**

a) For each \( n \), \( T_n \) is primitive recursive.

b) There is a primitive recursive function \( U \) such that if \( t \) is a computation tree, then \( U(t) \) is the output of the corresponding computation.

c) (Kleene’s Normal Form Theorem)

For every arity \( n \) and all \( e \) we have

\[
\phi_e(x_1, \ldots, x_n) = U(\mu t. T_n(e, x_1, \ldots, x_n, t)).
\]

**Proof**

It is only a) that requires a proof. The proof of a) is however easy, we construct the characteristic function of \( T_n \) by recursion on the last variable \( t \). Monotonicity of the sequence numbering is important here. We leave the tedious, but simple details for the reader.

**Corollary 3.2.25** For each number \( n \), the function

\[
f(e, x_1, \ldots, x_n) = \phi_e(x_1, \ldots, x_n)
\]

is computable.

**Remark 3.2.26** Corollary 3.2.25 is the analogue of the existence of a universal Turing machine, we can enumerate the computable functions in such a way that each computable function is uniformly computable in any of the numbers enumerating it. This universal function is partial. There is no universal function for the total computable functions, see Exercise 3.9.
The Recursion Theorem

The recursion theorem is one of the key insights in computability theory introduced by Kleene. In programming terms it says that we may define a set of procedures where we in the definition of each procedure refer to the other procedures in a circular way. The proof we give for the recursion theorem will be a kind of ‘white rabbit out of the hat’ argument based on the much more intuitive $S^n_m$-theorem. So let us first explain the $S^n_m$-theorem. Let $f$ be a computable function of several variables. Now, if we fix the value of some of the variables, we will get a computable function in the rest of the variables. The $S^n_m$-theorem tells us that the index for this new function can be obtained in a primitive recursive way from the index of the original function and the values of the fixed variables. We have to prove one technical lemma:

**Lemma 3.2.27** There is a primitive recursive function $\rho$ (depending on $n$) such that if

$$\phi^n_c(x_1, \ldots, x_n) = t$$

and

$$1 \leq i \leq n$$

then

$$\phi^{n-1}_{\rho(e, i, x_i)}(x_1, \ldots, x_{i-1}, x_{i+1}, \ldots, x_n) = t.$$ 

**Proof**

We define $\rho$ by induction on $e$, considering the cases i) - vi). We leave some of the easy details for the reader, see Exercise 3.11.

If $1 < i$ let $\rho(e, i, x_i) = e$, while $\rho(e, 1, x_1) = (3, x_1 + 1)$. The cases ii) and iii) are left for the reader.

If $e = (4, e', d_1, \ldots, d_n)$, we simply let

$$\rho(e, i, x_i) = (4, e', \rho(d_1, i, x_1), \ldots, \rho(d_n, i, x_i)).$$

Case v) splits into two subcases. If $1 < i$, this case is easy. For $i = 1$ we let

$$\rho(e, 1, 0) = d_1,$$

$$\rho(e, 1, x + 1) = (4, (3, x), \rho(e, 1, x), d_x, (2, 1), \ldots, (2, n))$$

where $(3, x)$ is the index for $f(\bar{x}) = x$ and $(2, i)$ is the index for the function selecting $x_i$ from $\bar{x}$. In this case the primitive recursion is replaced by an iterated composition, the depth of which is determined by the value of $x$.

Case vi) is again easy, and is left for the reader.

**Theorem 3.2.28 (The $S^n_m$-theorem)**

Let $n \geq 1, m \geq 1$. There is a primitive recursive function $S^n_m$ such that for all $c, x_1, \ldots, x_n, y_1, \ldots, y_m$

$$\phi^{n+m}_c(x_1, \ldots, x_n, y_1, \ldots, y_m) = \phi^{m}_{S^n_m(c, x_1, \ldots, x_n)}(y_1, \ldots, y_m)$$

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Proof
Let $\rho$ be as in Lemma 3.2.27.
Let $S^1_m(e, x) = \rho(e, 1, x)$
Let $S^k_{m+1}(e, x_1, \ldots, x_{k+1}) = \rho(S^k_m(e, x_1, \ldots, x_k), 1, x_{k+1})$. By an easy induction on $k$ we see that this construction works for all $m$.

The $S^m_n$-theorem is a handy tool in itself, and we will use it frequently stating that we can find an index for a computable function uniformly in some parameter. Now we will use the $S^m_n$-theorem to prove the surprisingly strong

**Theorem 3.2.29 (The Recursion Theorem)**
Let $f(e, \vec{x})$ be a partial, computable function.
Then there is an index $e_o$ such that for all $\vec{x}$
$$\phi_{e_o}(\vec{x}) \simeq f(e_0, \vec{x}).$$

**Proof**
Recall that by this equality we mean that either both sides are undefined or both sides are defined and equal. Let
$$g(e, \vec{x}) = f(S^1_n(e, e), \vec{x})$$
and let $\hat{g}$ be an index for $g$. Let $$e_0 = S^1_n(\hat{g}, \hat{g}).$$
Then
$$\phi_{e_0}(\vec{x}) = \phi_{S^1_n(\hat{g}, \hat{g})}(\vec{x}) = \phi_{\hat{g}}(\hat{g}, \vec{x}) = g(\hat{g}, \vec{x}) = f(S^1_n(\hat{g}, \hat{g}), \vec{x}) = f(e_0, \vec{x}).$$

**Remark 3.2.30** Readers familiar with the fixed point construction in untyped $\lambda$-calculus may recognize this proof as a close relative, and indeed it is essentially the same proof. The recursion theorem is a very powerful tool for constructing computable functions by self reference. In Chapter 4 we will use the recursion theorem to construct computable functions essentially by transfinite induction. Here we will give a completely different application, we will prove that there is no nontrivial set of partial computable functions such that the set of indices for functions in the class is computable.

**Theorem 3.2.31 (Rice)**
Let $A \subseteq \mathbb{N}$ be a computable set such that if $e \in A$ and $\phi_e = \phi_d$ then $d \in A$.
Then $A = \mathbb{N}$ or $A = \emptyset$.

**Proof**
Assume not, and let $a \in A$ and $b \notin A$.
Let $f(e, x) = \phi_b(x)$ if $e \in A$ and $f(e, x) = \phi_a(x)$ if $e \notin A$.
By the recursion theorem, let $e_0$ be such that for all $x$
$$f(e_0, x) = \phi_{e_0}(x).$$
If $e_0 \in A$, then $\phi_{e_0} = \phi_b$ so $e_0 \notin A$.
If $e_0 \notin A$, then $\phi_{e_0} = \phi_a$ so $e_0 \in A$.
This is a clear contradiction, and the theorem is proved.
Corollary 3.2.32 (Unsolvability of the halting problem)
\(\{ (e, x) : \phi_e(x) \downarrow \} \) is not computable.

Remark 3.2.33 Rice’s theorem is of course stronger than the unsolvability of
the Halting Problem, for which we need much less machinery.

3.2.3 Computably enumerable sets

Four equivalent definitions
An enumeration of a set \( X \) is an onto map \( F : \mathbb{N} \rightarrow X \). For subsets of \( \mathbb{N} \) we
may ask for computable enumerations of a set. A set permitting a computable
enumeration will be called computably enumerable or just c.e. The standard
terminology over many years has been recursively enumerable or just r.e., be-
cause the expression recursive was used by Kleene and many with him. We will
stick to the word computable and thus to the term computably enumerable.
Of course there is no enumeration of the empty set, but nevertheless we will
include the empty set as one of the c.e. sets.
In this section we will give some characterizations of the c.e. sets. One impor-
tant characterization will be as the semi-decidable sets. A computable set
will be decidable, we have an algorithm for deciding when an element is in the set
or not. In a semi-decidable set we will have an algorithm that verifies that an
element is in the set when it is, but when the element is not in the set, this
algorithm may never terminate. A typical semi-decidable set is the solving set
of the Halting Problem
\( \{ (e, x) : \phi_e(x) \downarrow \} \).
Another example is the set of theorems in first order number theory or any
nicely axiomatizable theory. The set of words in some general grammar will
form a third class of examples.

We will show that the semi-decidable subsets of \( \mathbb{N} \) will be exactly the c.e.
sets. A third characterization will be that the c.e. sets are exactly the sets of
projections of primitive recursive sets. In the literature this class is known as the
\( \Sigma^0_1 \)-sets. A fourth characterization will be as the ranges of partial, computable
functions.

This is enough talk, let us move to the definition:

Definition 3.2.34 Let \( A \subseteq \mathbb{N} \). We call \( A \) computably enumerable or just c.e.
if \( A = \emptyset \) or \( A \) is the range of a total computable function.

Theorem 3.2.35 Let \( A \subseteq \mathbb{N} \). Then the following are equivalent:

i) \( A \) is c.e.

ii) \( A \) is the range of a partial computable function.

iii) There is a primitive recursive set \( S \subseteq \mathbb{N}^2 \) such that

\[
A = \{ n : \exists m (n, m) \in S \}
\]
iv) There is a partial computable function with domain $A$.

Proof
Since the empty set satisfies all four properties, we will assume that $A \neq \emptyset$.

i) $\Rightarrow$ ii):
Trivial since we in this compendium consider the total functions as a subclass of the partial functions.

ii) $\Rightarrow$ iii):
Let $A$ be the range of $\phi_e$.

Then

$$n \in A \iff \exists y (T(e, \pi_1(y), \pi_2(y)) \land n = U(\pi_2(y)))$$

where $\pi_1$ and $\pi_2$ are the inverses of the pairing function $P$, $T$ is Kleene’s $T$-predicate and $U$ is the function selecting the value from a computation tree.

The matrix of this expression is primitive recursive.

iii) $\Rightarrow$ iv):
Let

$$n \in A \iff \exists m ((n, m) \in S).$$

where $S$ is primitive recursive. Then $A$ is the domain of the partial computable function

$$f(n) = \mu m. (n, m) \in S.$$  

iv) $\Rightarrow$ i):
Let $A$ be the domain of $\phi_e$ and let $a \in A$ (here we will use the assumption that $A$ is non-empty).

Let $f(y) = \pi_1(y)$ if $T(e, \pi_1(y), \pi_2(y))$, $f(y) = a$ otherwise. Then $f$ will be computable, and $A$ will be the range of $f$.

This ends the proof of the theorem.

Clearly characterizations ii) and iii) make sense for subsets of $\mathbb{N}^n$ as well for $n > 1$, and the equivalence will still hold. From now on we will talk about c.e. sets of any dimension. The relationship between c.e. subsets of $\mathbb{N}$ and $\mathbb{N}^n$ is given in Exercise 3.14.

The following lemma will rule our abilities to construct c.e. sets:

**Lemma 3.2.36** Let $A \subseteq \mathbb{N}^n$ and $B \subseteq \mathbb{N}^n$ be c.e. Then

a) $A \cap B$ and $A \cup B$ are both c.e.

b) If $n = m + k$ and both $m$ and $k$ are positive, then

$$\{ \bar{x} ; \exists \bar{y} (\bar{x}, \bar{y}) \in A \}$$

will be c.e., where $\bar{x}$ is a sequence of variables of length $m$ and $\bar{y}$ is of length $k$.

Moreover, every computable set is c.e., the inverse image or direct image of a c.e. set using a partial computable function will be c.e.
All these claims follow trivially from the definition or from one of the characterizations in Theorem 3.2.35.

Let us introduce another standard piece of notation:

**Definition 3.2.37** Let

\[ W_e = \{ n : \phi_e(n) \downarrow \} = \{ n : \exists t T(e, n, t) \}. \]

Let

\[ W_{e,m} = \{ n : \exists t < m T(e, n, t) \}. \]

We let \( \mathcal{K} \) be the diagonal set

\[ \mathcal{K} = \{ e : e \in W_e \}. \]

**Lemma 3.2.38**

(a) \( \{(e, n) : n \in W_e \} \) is c.e.

(b) \( \mathcal{K} \) is c.e.

(c) Each set \( W_{e,m} \) is finite.

(d) \( \{(e, n, m) : n \in W_{e,m} \} \) is primitive recursive.

All proofs are trivial.

**Selection with consequences**

From now on we will prove lemmas and theorems in the lowest relevant dimension, but clearly all results will hold in higher dimensions as well.

**Theorem 3.2.39 (The Selection Theorem)**

Let \( A \subseteq \mathbb{N}^2 \) be c.e. Then there is a partial computable function \( f \) such that

(i) \( f(n) \downarrow \iff \exists m (n, m) \in A. \)

(ii) If \( f(n) \downarrow \) then \( (n, f(n)) \in A. \)

**Proof**

This time we will give an intuitive proof. Let \( A \) be the projection of the primitive recursive set \( B \subseteq \mathbb{N}^3 \) (characterization iii). For each \( n \), search for the least \( m \) such that \( (n, \pi_1(m), \pi_2(m)) \in B \), and then let \( f(n) = \pi_1(m) \).

Intuitively we perform a parallel search for a witness to the fact that \( (n, m) \in A \) for some \( m \), and we choose the \( m \) with the least witness.

**Corollary 3.2.40** Let \( A \) and \( B \) be two c.e. sets. Then there are disjoint c.e. sets \( C \) and \( D \) with

\[ C \subseteq A, \ D \subseteq B \text{ and } A \cup B = C \cup D. \]
Proof
Let \( E = (A \times \{0\}) \cup (B \times \{1\}) \) and let \( f \) be a selection function for \( E \).
Let \( C = \{ n : f(n) = 0 \} \) and \( D = \{ n : f(n) = 1 \} \).
Then \( C \) and \( D \) will satisfy the properties of the corollary.

**Corollary 3.2.41** A set \( A \) is computable if and only if \( A \) and the complement of \( A \) are c.e.

One way is trivial, since the complement of a computable set is computable and all computable sets are c.e. So assume that \( A \) and its complement \( B \) are c.e. Let \( E \) be as in the proof of the corollary above, and \( f \) the selection function. Then \( f \) is the characteristic function of \( A \), so \( A \) is computable.

**Corollary 3.2.42** Let \( f : \mathbb{N} \rightarrow \mathbb{N} \) be a partial function. Then the following are equivalent:

i) \( f \) is computable.

ii) The graph of \( f \) is c.e.

**Proof**
If the graph of \( f \) is c.e., then \( f \) will be the selection function of its own graph, which is computable by the selection theorem. If on the other hand \( f \) is computable, then the graph of \( f \) will be the domain of the following function \( g(n, m) \):
Compute \( f(n) \) and see if the result equals \( m \).

**Computably inseparable c.e. sets**
In Exercise 3.16 we will see that two disjoint complements of c.e. sets can be separated by a computable set. This is an improvement of Corollary 3.2.42. Here we will show that a similar separation property does not hold for c.e. sets, and we will draw some consequences of this fact.

**Definition 3.2.43** Let \( A \) and \( B \) be two disjoint subsets of \( \mathbb{N} \). We say that \( A \) and \( B \) are computably separable if there is a computable set \( C \) such that \( A \subseteq C \) and \( B \cap C = \emptyset \). Otherwise \( A \) and \( B \) are computably inseparable.

**Theorem 3.2.44** There is a pair of computably inseparable c.e. sets.

**Proof**
Let \( A = \{ e : \phi_e(e) = 0 \} \) and \( B = \{ e : \phi_e(e) = 1 \} \).
Assume that \( C \) is a computable set that separates \( A \) and \( B \), and assume that \( e_0 \) is an index for the characteristic function of \( C \). Then, if \( e_0 \in C \) we have that \( \phi_{e_0}(e_0) = 1 \). Then \( e_0 \in B \) which is disjoint from \( C \).
Likewise, if \( e_0 \notin C \) we see that \( e_0 \in A \subseteq C \). In both cases we obtain a contradiction, so the existence of \( C \) is impossible.

Now this theorem has some interesting consequences concerning the difference between classical and constructive mathematics. We will end our general
introduction to the basics of computability theory discussing some of the consequences.

**Definition 3.2.45**  
a) A binary tree is a non-empty set $D$ of finite 0-1-sequences such that any initial segment of an element in $D$ is also in $D$. A binary tree is computable if the set of sequence numbers of the sequences in $D$ is computable.

b) An infinite branch in a binary tree $D$ is a function $f : \mathbb{N} \rightarrow \{0,1\}$ such that $(f(0), \ldots, f(n-1)) \in D$ for all $n$.

**Lemma 3.2.46 (König’s Lemma)**  
An infinite binary tree has an infinite branch.

The proof of König’s lemma is trivial, but has very little to do with computability theory. One constructs a branch by always extending the sequence in a direction where the tree is still infinite.

Well known theorems proved by similar arguments will be that a continuous function on a closed bounded interval will obtain its maximum and that any consistent first order theory has a complete extension.

**Remark 3.2.47** The version of König’s Lemma given above, restricting ourselves to binary trees, is often called Weak König’s lemma. The full König’s lemma then refers to infinite trees with finite branching, not just binary branching. There are results showing that Weak König’s lemma is equivalent, relative to some very weak formal theory, to the mentioned theorems from analysis and topology.

We will show a failure of a constructive version of König’s lemma:

**Lemma 3.2.48** There is an infinite, computable binary tree without a computable, infinite branch.

Proof  
Let $A$ and $B$ be two computably inseparable c.e. sets and let $\{A_n\}_{n \in \mathbb{N}}$ and $\{B_n\}_{n \in \mathbb{N}}$ be primitive recursive sequences of sets $A_n$ and $B_n$, contained in $\{0, \ldots, n-1\}$, with $A = \bigcup_{n \in \mathbb{N}} A_n$ and $B = \bigcup_{n \in \mathbb{N}} B_n$.

If $\sigma$ is a 0-1-sequence of length $n$, we let $\sigma \in D$ if for all $i < n$: $i \in A_n \Rightarrow \sigma(i) = 1$ and $i \in B_n \Rightarrow \sigma(i) = 0$. $\sigma$ will be a characteristic function on $\{0, \ldots, n-1\}$ separating $A_n$ and $B_n$.

Now the characteristic function of any set separating $A$ and $B$ will be an infinite branch in $D$. Since $A$ and $B$ are disjoint, $D$ will be an infinite, binary tree. Moreover, $D$ will be primitive recursive. A computable, infinite branch will on the other hand be the characteristic function of a computable set separating $A$ and $B$, something assumed not to exist. This ends the proof.

In Exercises 3.17 and 3.18 we will give an application of this lemma to propositional calculus and to constructive analysis.
3.2.4 Crosses and naughts

Once upon the time there was a popular game known as “Crosses and naughts”. This game requires two players and a potentially large piece of paper split into squares. Player no. 1 writes a cross in one of the squares, player no. 2 writes a naught in one of the empty squares, player no. 1 writes a cross in one of the empty squares and so on. The first player who fills a row of five symbols, up-down, sideways or along a diagonal has won the game.

In order to turn this into a mathematical problem, we make two assumptions:

- The sheet of paper is infinite in all directions.
- Each player must write a symbol within a prefixed distance of the game played so far at each move, e.g. in a square sharing at least one point with a square already filled.

We will see that it is impossible for player 2 to have a winning strategy. If player 2 had a winning strategy, player 1 could pretend that she/he was player 2 and base her/his game on that strategy and on an imagined first move from the real player 2. If player 2 later makes this imagined move, player 1 just base her/his further game on a new imagined move. This will enable player 1 to win.

This leaves us with two possibilities:

1. Player 1 will have a winning strategy.
2. Player 2 has a strategy for dragging the game into infinite time.

It follows by weak König’s lemma that one of the two will hold.

**Conjecture 3.2.49** Player 2 has a strategy for a never-ending game, but there is no computable strategy to this effect.

It is a consequence of König’s lemma that if Player 1 has a winning strategy, there is an upper limit to how man moves Player 1 will need in order to win the game. Then we may even abandon the assumption that the sheet of paper is infinite. To the author, this seems rather unlikely. This supports one part of the conjecture. The other part is not really supported by a belief, but by the appealing possibility that you might beat any computer in this game, if you get hold of how it is programmed, but you may not be able to beat an improvising human being.

3.3 Degrees of Unsolvability

3.3.1 m-reducibility

Discussion

Having introduced the main concepts of computability theory there are several options. One option will be to give a further analysis of the computable functions.
and subclasses of them. This will lead us to complexity theory or to subrecursive hierarchies (see Section 3.6). One key concept then will be that one set or function can be reduced to another in some simple way.

We will not consider such a fragmentation of the computable functions into interesting subclasses yet. Instead we will ask for reductions between possible non-computable functions and sets, using the complexity of computability itself for defining the notion of reduction. The connection between this approach and the ones mentioned above is that both the formation of interesting concepts and the methodology of those approaches to computability theory are based on the experience gained from investigating computable reductions in general. Thus, in an introductory course, where one should learn the computability theorists way of thinking, this section is basic.

**Definition 3.3.1** Let $A$ and $B$ be two subsets of $\mathbb{N}$.
We say that $A$ is \textit{m-reducible to} $B$, $A <_m B$, if there is a total computable function $f$ such that for all $n$:

$$n \in A \iff f(n) \in B.$$ 

We read this as ‘many-one-reducible’, since $f$ may be a many-one function. If we insist on $f$ being injective, we will get 1-reducibility.

There is no way we can reduce $\mathbb{N}$ to $\emptyset$ and vice versa, and we will exclude those trivial sets from our further discussion. We then get the observations:

**Lemma 3.3.2**

a) If $A$ is computable and $B \neq \mathbb{N}, \emptyset$, then $A <_m B$.

b) $A <_m B \iff (\mathbb{N} \setminus A) <_m (\mathbb{N} \setminus B)$.

c) $<_m$ is transitive.

The proofs are easy and are left for the reader.

**Definition 3.3.3** We call two sets $A$ and $B$ \textit{m-equivalent} if $A <_m B$ and $B <_m A$. We then write $A \equiv_m B$.

Clearly $\equiv_m$ will be an equivalence relation with a partial ordering inherited from $<_m$. We call the set of equivalence classes with this induced ordering the \textit{m-degrees}.

**Remark 3.3.4** Let us briefly discuss the distinction between a \textit{complexity class} and a \textit{degree}. A complexity class will consist of all functions and/or sets that may be computed or decided with the aid of a prefixed amount of computing resources. Resources may be the time or space available (normally as functions of the input), but to us it may also be the possible service of a non-computable oracle. This will be better explained later. A degree, on the other hand, will be some mathematical object measuring the complexity of a function or a set. It has turned out that the most useful object for this purpose simply is the

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class of functions of the same complexity as the given object. It is like the old Frege semantics where 17 is interpreted as the set of all sets with 17 elements, and in general a property is interpreted as the set of all objects sharing that property. It remains to decide what we mean by ‘having the same complexity’, and we will consider two (out of many in the literature) possible choices. Being \( m \)-equivalent is one, and later we will learn about Turing-equivalent functions.

The \( m \)-degrees have some properties easy to establish:

**Lemma 3.3.5**  
\( a) \) Let \( X \) be a finite set of \( m \)-degrees. Then \( X \) has a least upper bound.  
\( b) \) Let \( X \) be a countable set of \( m \)-degrees. Then \( X \) has an upper bound.

**Proof**  
The computable sets will form the minimal \( m \)-degree. Thus the empty set of \( m \)-degrees has a least upper bound. In order to prove the rest of \( a) \) it is sufficient to show that \( \{ A, B \} \) will have a least upper bound. Let  
\[
A \oplus B = \{ 2n \mid n \in A \} \cup \{ 2n + 1 \mid n \in B \}.
\]
This will define the least upper bound, see Problem 3.19.  
In order to prove \( b) \), Let \( \{ A_n \mid n \in \mathbb{N} \} \) be a countable family of sets. Let  
\[
A = \Sigma_{n \in \mathbb{N}} A_n = \{ (n,m) \mid m \in A_n \}.
\]
Then the \( m \)-degree of \( A \) will bound the \( m \)-degrees of \( A_n \) for each \( n \in \mathbb{N} \).

We will invest most of our efforts in analyzing the Turing degrees below. We will however prove one negative result, the \( m \)-degrees do not represent a linear stratification of all nontrivial sets into complexity classes.

**Lemma 3.3.6**  
There are two nontrivial sets \( A \) and \( B \) such that \( A \not\leq_m B \) and \( B \not\leq_m A \).

**Proof**  
Let \( A \) be c.e. but not computable. Let \( B = \mathbb{N} \setminus A \).  
If \( C \leq_m A \), \( C \) will be c.e., so \( B \) is not \( m \)-reducible to \( A \). It follows from Lemma 3.3.2 that \( A \) is not \( m \)-reducible to \( B \).

**An \( m \)-complete c.e. set**  
We proved that the computable inverse of a c.e. set is c.e. A reformulation of this will be

**Lemma 3.3.7**  
Let \( A <_m B \). If \( B \) is c.e., then \( A \) is c.e.

Thus the c.e. sets form an initial segment of all sets pre-ordered by \( m \)-reductions. First we will show that this set has a maximal element:
Lemma 3.3.8 Let $K = \{ e : \phi_e(e) \downarrow \}$. Then any other c.e. set is $m$-reducible to $K$.

Proof
Let $A = \{ d : \phi_d(d) \downarrow \}$. Adding a dummy variable we can without loss of generality assume that $A = \{ d : \phi_e(d, x) \downarrow \}$ where the outcome of $\phi_e(d, x)$ is independent of $x$. Then $A <_m K$ by

$$A = \{ d ; S^1(e, d) \in K \}.$$

A natural question is now if there are c.e. sets that are not $m$-equivalent to $K$ or the computable sets; are there more $m$-degrees among the c.e. sets? This was answered by Emil Post, when he classified a whole class of c.e. sets 'in between'.

Simple sets

Definition 3.3.9 A c.e. set $A$ is simple if the complement of $A$ is infinite, but does not contain any infinite c.e. sets.

A simple set cannot be computable (why?). We will prove that there exist simple sets, and that $K$ cannot be reduced to any simple sets.

Lemma 3.3.10 There exists a simple set.

Proof
Let $B = \{ (e, x) : 2e < x \land \phi_e(x) \downarrow \}$. Let $g$ be a computable selection function for $B$, i.e. $g(e) \downarrow$ when $(e, x) \in B$ for some $x$, and then $g(e)$ selects one such $x$. Let $A$ be the image of $g$. Then $A$ is c.e.

Since $g(e) > 2e$ when defined, the complement of $A$ will be infinite. This is because each set $\{ 0, \ldots, 2e \}$ will contain at most $e$ elements from $A$. If $W_e$ is infinite, $W_e$ will contain a number $> 2e$, and $g(e)$ will be defined. Then $A \cap W_e \neq \emptyset$, so $W_e$ is not contained in the complement of $A$. This shows that $A$ is simple.

Lemma 3.3.11 Let $K <_m A$ where $A$ is c.e. Then the complement of $A$ contains an infinite c.e. set.

Proof
Let $f$ be total and computable such that

$$e \in K \Leftrightarrow f(e) \in A.$$

We will construct a computable sequence $\{ x_n \}_{n \in \mathbb{N}}$ of distinct numbers outside $A$, and use the recursion theorem to glue the whole construction together. The induction start will be the empty sequence.
Assume that \( B_n = \{ x_0, \ldots, x_{n-1} \} \) has been constructed such that \( B_n \) is disjoint from \( A \).

Let \( \rho(n) \) be such that \( W_{\rho(n)} = \{ e : f(e) \in B_n \} \). By the uniformity of the construction, \( \rho \) will be computable. We will let \( x_n = f(\rho(n)) \).

If \( x_n \in B_n \) we have that \( x_n \notin A \), so \( \rho(n) \notin \mathcal{K} \), and \( \rho(n) \notin W_{\rho(n)} \). This contradicts the assumption that \( x_n \in B_n \) and the definition of \( W_{\rho(n)} \). Thus \( x_n \notin B_n \).

On the other hand, if \( x_n \in A \), then \( \rho(n) \in \mathcal{K} \), \( \rho(n) \in W_{\rho(n)} \) and by construction, \( x_n \in B_n \), which we agreed was impossible. Thus \( x_n \) is a new element outside \( A \). We can then continue the process and in the end produce an infinite c.e. set outside \( A \).

**Corollary 3.3.12** There is a non-computable c.e. set \( A \) such that \( \mathcal{K} \) is not \( m \)-reducible to \( A \).

### 3.3.2 Turing degrees

**Relativised computations**

In the previous section we considered a notion of reducibility between sets based on the idea ‘\( A \) is simpler than \( B \) if we can gain information about \( A \) by making one computable call about membership in \( B \)’. Now, if we imagine that we are given the superhuman ability to decide membership in \( B \), we must also permit ourselves to ask more intricate questions about \( B \) in order to decide membership in \( A \), and still claim that we use our superhuman power to make \( A \) decidable. There are several notions of reductions based on the idea that we may be allowed any computable set of questions about membership in \( B \) and draw the conclusions from the answers. We will not consider any of those intermediate notions of reduction in this introductory course, but take the most drastic approach: We extend the definition of the computable functions by adding some extra functions as initial ones. This is technically described in the following definition:

**Definition 3.3.13** Let \( f_1, \ldots, f_k \) be a finite list of partial functions \( f_i : \mathbb{N} \rightarrow \mathbb{N} \).

a) We extend the definition of \( \phi_e(\vec{x}) \) to a definition of \( \phi_e^{f_1, \ldots, f_k}(\vec{x}) \)

by adding a new clause

\[ \text{ vii) If } e = (7, i, k) \text{ then } \phi_e^{f_1, \ldots, f_k}(x, \vec{x}) = f_i(x). \text{ If } f_i(x) \text{ is undefined, then this computation does not terminate.} \]

b) The computation tree will in this case be \( (7, i, k, x, f_i(x)) \).

We call these algorithms relativized algorithms meaning that the concept of computation has been relativized to \( f_1, \ldots, f_k \). One important aspect of this definition is the finite use principle. Even if we in reality accept functions as inputs in algorithms as well, no terminating algorithm will use more than a finite amount of information from the functions involved.
Lemma 3.3.14 Assume that $\phi^{f_1,\ldots,f_k}_{e}(\vec{x}) = z$. Then there are finite partial subfunctions $f'_1,\ldots,f'_k$ of $f_1,\ldots,f_k$ resp. such that

$$
\phi^{f'_1,\ldots,f'_k}_{e}(\vec{x}) = z
$$

with the same computation tree.

The proof is by a tedious but simple induction on $e$, using the fact that a sequence number is larger than all parts of the sequence.

We have given and used an enumeration of all ordered sequences of natural numbers. In the same fashion we can construct an enumeration $\{\xi_i\}_{i \in \mathbb{N}}$ of all finite partial functions such that the following relations will be primitive recursive:

- $\xi_i(x) \downarrow$
- $\xi_i(x) = y$
- $\text{dom}(\xi_i) \subseteq \{0,\ldots,n\}$.

The reader is free to add any other important properties, as long as they are correct.

Example 3.3.15 In order to avoid that the reader will focus on inessential details, we offer the construction of $\xi_i$ in the guise of an example.

Let $p_0, p_1, p_2, \ldots$ be an enumeration of the prime numbers, starting with $p_0 = 2$, $p_1 = 3$ and so on.

Define $\xi_i$ by

$$
\xi_i(n) = m \text{ if there are exactly } m + 1 \text{ factors of } p_n \text{ in } i.
$$

$$
\xi_i \uparrow \text{ if } p_n \text{ is not a factor in } i.
$$

For example, if $i = 2^3 \cdot 7^4 \cdot 11$, then $\xi_i(0) = 2, \xi_i(3) = 3$ and $\xi_i(4) = 0$. For all other inputs $x$, $\xi_i(x)$ will be undefined.

The properties above are sufficient to prove the extended Kleene T-predicate:

Lemma 3.3.16 For each $n$ and $k$ the following relation is primitive recursive:

$$
T_{n,k}(e,\vec{x},i_1,\ldots,i_k,t) \iff t \text{ is the computation tree of } \phi^{\xi_{i_1},\ldots,\xi_{i_k}}_{e}(\vec{x}).
$$

In order to save notation we will from now on mainly consider computations relativized to one function. Using the coding

$$
\langle f_0,\ldots,f_{k-1} \rangle(x) = f_{\pi_1(x)(\mod k)}(\pi_2(x))
$$

we see that there is no harm in doing this.

Lemma 3.3.17 There is a primitive recursive function $c$ such that for any partial functions $f$, $g$ and $h$, if for all $x$, $f(x) = \phi^g_e(x)$ and for all $x$, $g(x) = \phi^h_d(x)$, then for all $x$, $f(x) = \phi^h_{c(e,d)}(x)$. 

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Proof

c(e, d) is defined by primitive recursion on e, dividing the construction into cases i) - xi).

For cases i) - iii), we let \( c(e, d) = e \).

For the cases iv) - viii) we just let \( c(e, d) \) commute with the construction of \( e \) from its subindices, i.e., if the case is \( e = expr(e_1, \ldots, e_n) \) then \( c(e, d) = expr(c(e_1, d), \ldots, c(e_n, d)) \), where \( expr \) can be any relevant expression.

If \( e = \langle 9, 1, 1 \rangle \) we have \( \phi^{g}_{e}(x, \vec{x}) = g(x) = \phi^{h}_{d}(x) \), so we let \( c(e, d) = d' \) where \( \phi^{h}_{d'}(x, \vec{x}) = \phi^{d}_{e}(x) \). \( d' \) is primitive recursive in \( d \).

**Remark 3.3.18** There is an alternative proof of this lemma. By the finite use principle we see that if \( g \) is computable in \( h \) and \( f \) is computable in \( g \), then the graph of \( f \) will be c.e. relative to \( h \). Everything will be uniform, so we may extract \( c \) from this proof as well.

We can use the concepts introduced here to talk about computable functionals, and not just about computable functions.

**Definition 3.3.19** Let \( F : \mathbb{N}^{\mathbb{N}} \to \mathbb{N} \), a functional of type 2.
We call \( F \) computable if there is an index \( e \) such that for all total \( f : \mathbb{N} \to \mathbb{N} \) we have
\[
F(f) = \phi^{f}_{e}(0).
\]
All computable functionals of type 2 will be continuous with respect to the canonical topologies. Kleene extended the notion of relativized computation to cover all functionals of any finite type as possible inputs. This will be discussed briefly in Chapter 4.

**The second recursion theorem**

We are now in the position to state and prove the second recursion theorem:

**Theorem 3.3.20** Let \( F(f, x) = \phi^{f}_{e}(x) \) be a computable functional.
Then there is a computable \( f \) such that for all \( x \),
\[
f(x) = F(f, x).
\]
Moreover, \( f \) will be the least such function in the sense that if
\[
\forall x(g(x) = F(g, x)),
\]
then \( f \subseteq g \) as graphs.

**Proof**
Let \( f_0 \) be the everywhere undefined function, and by recursion let \( f_{n+1}(x) = F(f_n, x) \).
Since \( f_0 \subseteq f_1 \) we see by induction on \( n \) that \( f_n \subseteq f_{n+1} \) for all \( n \).
Let \( f = \bigcup\{f_n : n \in \mathbb{N} \} \). By the finite use principle we see that
\[
F(f, x) = y \iff \exists n F(f_n, x) = y \iff \exists n f_{n+1}(x) = y \iff f(x) = y.
\]
Clearly \( f \) is the minimal solution to the equation, the minimal fixed point of \( F \).
It remains to show that \( f \) is computable, i.e. that the graph of \( f \) is c.e.
We will use the effective enumeration of all finite partial functions and the effective coding of sequences.
Let \( \xi \) with indices vary over all finite partial functions.
Then we can rewrite the following characterization of the graph of \( f \) in a \( \Sigma^0_1 \)-form:

\[
\forall i < n \exists \exists (\xi_1, \ldots, \xi_n) (\forall z \xi_0(z)) \\
\land \forall i < n (\forall z \in \text{dom}(\xi_{i+1}) (\xi_{i+1}(z) = F(\xi_i, z)) \land \xi_n(x) = y).
\]

**Remark 3.3.21** The existence of a computable fixed point for \( F \) follows by the recursion theorem.
It is possible to show that the index obtained by the proof of the recursion theorem will be an index for the least fixed point, see Exercise 3.20.

**Turing Degrees**

We will now restrict ourselves to total functions.

**Definition 3.3.22** Let \( f \) and \( g \) be two functions.
We say that \( f \) is computable in \( g \) if there is an index \( e \) such that for all \( x \),

\[
f(x) = \phi^g_e(x).
\]

We write \( f <_T g \). We will then say that \( f \) is Turing reducible to \( g \).

The key properties of Turing reducibility is given by

**Lemma 3.3.23**

a) \( <_T \) is a transitive relation.

b) If \( f \) is computable and \( g \) is any function, then \( f <_T g \).

c) If \( f \) and \( g \) are functions, there is a function \( h \) such that for any other function \( h' \):

\[
f <_T h \text{ and } g <_T h. \\
\text{If } f <_T h' \text{ and } g <_T h' \text{ then } h <_T h'.
\]

**Proof**

a) is a consequence of Lemma 3.3.17, b) is trivial, and to prove c), let \( h(2n) = f(n) \) and \( h(2n + 1) = g(n) \).

**Definition 3.3.24** Let \( f \) and \( g \) be two functions. \( f \) and \( g \) are Turing equivalent, in symbols \( f \equiv_T g \), if \( f <_T g \) and \( g <_T f \).

\( \equiv_T \) will be an equivalence relation. The equivalence classes will be called Turing degrees or Degrees of unsolvability. We will simply call them degrees. We will let bold-face low case letters early in the alphabet , \( a, b \), etc. denote degrees.
The set of degrees has a canonical ordering \( < \) inherited from \( <_T \).

We can summarize what we have observed so far in
Lemma 3.3.25 The ordered set of degrees is an upper semi-lattice with a least element such that every countable set is bounded, and every initial segment is countable.

We leave the verifications for the reader. We will now see that there is no maximal degree:

Lemma 3.3.26 Let $a$ be a degree. Then there is a degree $b$ such that $a < b$.

Proof
Let $f \in a$. Let $g(x) = \phi^f_{\pi_1(x)}(\pi_2(x)) + 1$ if $\phi^f_{\pi_1(x)}(\pi_2(x)) \downarrow$, otherwise $g(x) = 0$. The proof of the unsolvability of the halting problem can be relativized to $f$ so $g$ is not computable in $f$. On the other hand, clearly $f$ is computable in $g$.

The $g$ constructed in the proof above is called $f'$, the jump of $f$. The jump operator is indeed a degree-operator, see Exercise 3.21.

We have shown that there are incomparable $m$-degrees. The same method can be used to show that there are incomparable Turing degrees, see Exercise 3.22. We will prove a stronger result, showing that no strictly increasing sequence of degrees will have a least upper bound.

Theorem 3.3.27 Let $\{a_i\}_{i \in \mathbb{N}}$ be a strictly increasing sequence of degrees. Then there are two degrees $b$ and $c$ that are both upper bounds for the sequence, such that for any degree $d$, if $d < b$ and $d < c$, then $d < a_i$ for some $i \in \mathbb{N}$.

A pair $b, c$ as above is called a perfect pair for the sequence. The degrees in a perfect pair for a sequence will be incomparable. Further, the existence of a perfect pair shows that the sequence will have no least upper bound.

Proof
Let $\{f_i\}_{i \in \mathbb{N}}$ be a sequence of total functions of increasing Turing degrees. We will construct the functions $g$ and $h$ as the limits of approximations $g_x$ and $h_x$, where we in the construction of $g_{x+1}$ and $h_{x+1}$ want to ensure that if $e_1 = \pi_1(x)$ and $e_2 = \pi_2(x)$ and $\phi^g_{e_1} = \phi^h_{e_2}$ are total, then $\phi^g_{e_1}$ is computable in $f_x$. In order to simplify the notation, we let $g$ and $h$ be defined on $\mathbb{N}^2$, but this will alter nothing.

In the construction we will preserve the following properties:

1. If $i < x$, then $g_x(i, n)$ and $h_x(i, n)$ are defined for all $n$.
2. If $i < x$, then $g_x(i, n) = f_i(n)$ for all but finitely many $n$.
3. If $i < x$, then $h_x(i, n) = f_i(n)$ for all but finitely many $n$.
4. $g_x(i, n)$ is defined only for finitely many $(i, n)$ with $x \leq i$.
5. $h_x(i, n)$ is defined only for finitely many $(i, n)$ with $x \leq i$.

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This will ensure that $g_x$ and $h_x$ are equivalent to $f_{x-1}$ for $x \geq 0$.

Let $g_0$ and $h_0$ both be the empty function.

Now, let $x \geq 0$ and assume that $g_x$ and $h_x$ are defined satisfying 1. - 5. above.

Let $e_1 = \pi_1(x)$ and $e_2 = \pi_2(x)$. What we will do next will depend on the answer to the following question:

Can we find an $n$ and finite extensions $g'$ of $g_x$ and $h'$ of $h_x$ such that $\phi_{e_1}^{g'}(n) \neq \phi_{e_2}^{h'}(n)$ and both are defined?

(By a finite extension we mean that we ad a finite number of elements to the domain.) If the answer is ‘no’, we extend $g_x$ to $g_{x+1}$ by letting $g_{x+1}(x, n) = f_x(n)$ whenever this is not in conflict with the construction of $g_x$ (a conflict that can appear at at most finitely many places), and we construct $h_{x+1}$ from $h_x$ and $f_x$ in the same way.

If the answer is ‘yes’, we first choose two such finite extensions, and then we construct $g_{x+1}$ and $h_{x+1}$ from these extensions as above. This ends our construction.

We let $g(x, n) = g_{x+1}(x, n)$ and $h(x, n) = h_{x+1}(x, n)$. In the construction we have tried as hard as possible to avoid that $\phi_e^{g'} = \phi_d^{h'}$. The point is that we have tried so hard that if they after all turn out to be equal, they will both be computable in one of the $f_i$’s.

Claim 1

$f_i$ is computable in both $g$ and $h$.

Proof

We have that $f_i(n) = g(i, n)$ except for finitely many $i$, so $f_i$ is computable in $g$. The same argument holds for $h$.

Claim 2

For $x \geq 0$ we have that $g_x$ and $h_x$ both are computable in $f_{x-1}$

Proof

This is a trivial consequence of properties 1. - 5.

Claim 3

If $\phi_e^g = \phi_d^h$ and both are total, then $\phi_e^g$ is computable in $f_x$ for some $x$.

Proof

Let $x = P(e, d)$ where $P$ is the pairing from Definition 3.2.13. Then in the construction of $g_{x+1}$ and $h_{x+1}$ we ask for a $y$ and finite extensions $g'$ and $h'$ of $g_x$ and $h_x$ such that $\phi_e^{g'}(y) \neq \phi_d^{h'}(y)$. If we had found some, we would let $g$ and $h$ be further extensions of one such pair of finite extensions, and then we would have preserved that $\phi_e^g(y) \neq \phi_d^h(y)$, contradicting our assumption on $e$ and $d$.

On the other hand, for every $y$ we can, by the assumption and the finite use principle, find finite extensions such that $\phi_e^g(y)$ and $\phi_d^h(y)$. The point is that all these values must be equal, otherwise we could have found two extensions giving different values. Thus we can give the following algorithm for computing $\phi_e^g(y)$ from $g_x$ which again is computable in $f_{x-1}$: Search for any finite extension (by searching through the finite partial functions consistent with $g_x$) $g'$ of $g_x$ such that $\phi_e^{g'}(y)$. The value we obtain will be the correct value.

This ends the proof of our theorem.
3.4 A minimal degree

3.4.1 Trees

In the constructions of degrees we have performed so far, we have been using brute force. For instance, when we want to construct a minimal pair, we start with three properties:

1. $f$ is not computable.
2. $g$ is not computable.
3. If $h$ is computable in both $f$ and $g$, then $h$ is computable.

We then split these properties into infinite lists of requirements, which we try to satisfy during the construction:

- $R_{1,e}$: If $\phi_e$ is total, then $f \neq \phi_e$.
- $R_{2,e}$: If $\phi_e$ is total, then $g \neq \phi_e$.
- $R_{3,e,d}$: If $\phi^f_e = \phi^g_d$ is total, then $\phi^f_e$ is computable.

Now, for any of these requirements and any pair $\sigma$ and $\tau$ of finite sequences there will be finite extensions $\sigma'$ and $\tau'$ such that any further total extension $f$ and $g$ of $\sigma'$ and $\tau'$ resp. will satisfy the requirement. Thus by a step-by-step construction we can construct $f$ and $g$ via finite approximations satisfying one requirement at the time. The reader is invited to work out the full proof, see Exercise 3.23.

We will now face a problem which we cannot solve by this simple method. We will show that there is a non-computable function $f$ such that there is no function of complexity strictly between $f$ and the computable functions. Again we will set up the relevant properties of $f$, and fragmentize them into a sequence of requirements we want to satisfy during the construction. The problem will be that we cannot ensure that these requirements are satisfied by considering just a finite approximation of $f$. Instead we will use trees, and we will satisfy the various requirements by insisting that $f$ is a branch in a given binary tree. Before we can go into details with the argument, we will reconsider our definition of a binary tree, give a formulation that will be handy for this particular application. We will not distinguish between finite sequences and sequence numbers here, but whenever we say that a function defined from a set of finite sequences to the set of finite sequences is computable, we mean that the corresponding function on sequence numbers is computable.

**Definition 3.4.1** Let $D$ be the set of finite 0-1-sequences.

a) If $\sigma \in D$ we let $\sigma * 0$ and $\sigma * 1$ be $\sigma$ extended by 0 or 1 resp. We extend this to the concatenation $\sigma * \tau$ in the canonical way.

b) If $\sigma$ and $\tau$ are two sequences, and $i < lh(\sigma)$, $i < lh(\tau)$ and $\sigma(i) \neq \tau(i)$, we say that $\sigma$ and $\tau$ are inconsistent.
c) $f : D \to D$ is monotone if $f(\sigma)$ is a proper subsequence of $f(\tau)$ whenever $\sigma$ is a proper subsequence of $\tau$.

d) A tree is a monotone function $T : D \to D$ mapping inconsistent sequences to inconsistent sequences.

e) If $S$ and $T$ are trees, then $S$ is a subtree of $T$ if there is a tree $T'$ such that $S = T \circ T'$.

f) If $T$ is a tree and $f : \mathbb{N} \to \{0, 1\}$, then $f$ is a branch in $T$ if for every $n$ there is a sequence $\sigma$ of length $n$ such that $T(\sigma)$ is an initial segment of $f$.

Remark 3.4.2 These trees will sometimes be called perfect trees, the set of branches will form a perfect subset of the set $\{0, 1\}^\mathbb{N}$ in the topological sense, i.e. a set that is closed and without isolated points.

Remark 3.4.4 Using topology we might just say that the intersection of a decreasing sequence of nonempty compact sets is nonempty, in order to prove this lemma.

Lemma 3.4.3 Let $\{T_n\}_{n \in \mathbb{N}}$ be a sequence of trees such that $T_{n+1}$ is a subtree of $T_n$ for all $n$. Then there is a function $f$ that is a branch in all trees $T_n$.

Proof Consider the set $X$ of $\sigma$ such that $\sigma \in \text{Set}(T_n)$ for all $n$. $X$ will be a binary tree. The empty sequence is in $X$. If $\sigma \in X$, then for all $n$, $\sigma * 0 \in \text{Set}(T_n)$ or $\sigma * 1 \in \text{Set}(T_n)$. At least one of these has to hold for infinitely many $n$ and, since we are dealing with subtrees, for all $n$. Thus $X$ is not a finite tree and by König’s Lemma has an infinite branch, which will be a common branch for all $T_n$’s.

Remark 3.4.5 Using topology we might just say that the intersection of a decreasing sequence of nonempty compact sets is nonempty, in order to prove this lemma.

We will now see that certain computable trees can be used to meet natural requirements. As a first case, let us prove:

Lemma 3.4.5 Let $T$ be a computable tree and assume that $\phi_e$ is total. Then there is a computable subtree $S$ of $T$ such that $\phi_e$ is not a branch in $S$.

Proof If $\phi_e$ is not a branch in $T$, we can use $S = T$. If $\phi_e$ is a branch in $T$, one of $T(0)$ and $T(1)$ will be inconsistent with $\phi_e$, since they are inconsistent themselves. (here 0 is the sequence of length 1 with entry 0). Assume that $\phi_e$ is inconsistent with $T(0)$.

Let $S(\sigma) = T(0 * \sigma)$. Then $S$ is a subtree as desired. The other case is essentially the same.
3.4.2 Collecting Trees

Now we will describe a property on computable trees that will ensure that if \( f \) is a branch in the tree and \( \phi^e_c \) is total, then \( \phi^e_c \) is computable.

**Definition 3.4.6** Let \( T \) be a computable tree, \( e \) an index. \( T \) is \( e \)-collecting if for all finite sequences \( \sigma, \tau \) and all \( x \in \mathbb{N} \), if \( \phi^T(\sigma)(x) \downarrow \) and \( \phi^T(\tau)(x) \downarrow \), then
\[
\phi^T(\sigma)(x) = \phi^T(\tau)(x).
\]

**Lemma 3.4.7** Let \( T \) be a computable \( e \)-collecting tree and let \( f \) be a branch in \( T \). If \( \phi^e_c \) is total, then \( \phi^e_c \) is computable.

**Proof**
We will give an algorithm for computing \( \phi^e_c(x) \) from \( x \).
Since \( \phi^e_c(x) \downarrow \), there will be a 0-1-sequence \( \sigma \) such that \( \phi^T(\sigma)(x) \downarrow \), and since \( T \) is \( e \)-collecting, the value \( \phi^T(\sigma)(x) \) will be independent of the choice of \( \sigma \). Thus our algorithm will be:
Search for a finite sequence \( \sigma \) such that \( \phi^T(\sigma)(x) \downarrow \) and let the answer be the output of our algorithm.

**Remark 3.4.8** There is more information to be gained from this proof. We see that the function \( \phi^e_c \) itself is independent of \( f \) as long as it is total and \( f \) is a branch in an \( e \)-collecting tree.

3.4.3 Splitting Trees

We will now find a criterion that will ensure that \( f \) is computable in \( \phi^e_c \) whenever \( f \) is a branch in a computable tree and \( \phi^e_c \) is total:

**Definition 3.4.9** Let \( T \) be a computable tree, and let \( e \) be an index.
We call \( T \) \( e \)-splitting if for all finite 0-1-sequences \( \sigma, \tau \), if \( \sigma \) and \( \tau \) are inconsistent, then there exists a number \( x \) such that
\[
\phi^T(\sigma)(x) \downarrow, \phi^T(\tau)(x) \downarrow \text{ with } \phi^T(\sigma)(x) = \phi^T(\tau)(x).
\]

**Lemma 3.4.10** Let \( T \) be an \( e \)-splitting computable tree, let \( f \) be a branch in \( T \) and assume that \( \phi^e_c \) is total. Then \( f \) is computable in \( \phi^e_c \).

**Proof**
We will compute an infinite 0-1-sequence \( \{k_i\}_{i \in \mathbb{N}} \) from \( \phi^e_c \) such that \( T(\sigma_n) \) is an initial segment of \( f \) for all \( n \), where \( \sigma_n = (k_0, \ldots, k_{n-1}) \). The empty sequence \( \sigma_0 \) of course satisfies this. Assume that \( \sigma_n \) is constructed. Then one of \( T(\sigma_n \ast 0) \) and \( T(\sigma_n \ast 1) \) will be an initial segment of \( f \). We will just have to determine which one. Now \( \phi^T(\sigma_n \ast 1) \) and \( \phi^T(\sigma_n \ast 1) \) will be inconsistent, so exactly one of them will be inconsistent with \( \phi^e_c \). We can then use \( \phi^e_c \) to find the inconsistent one, which means that we can decide which direction is along \( f \) and which is not. This provides us with the induction step, and we can move on. This ends the proof of the lemma.
Remark 3.4.11 In the case of $T$ being an $e$-splitting tree, we see that $\phi^e_T$ is a one-to-one-function of the branch $f$, and what we just argued for is that we can compute the inverse.

3.4.4 A minimal degree

Using Lemmas 3.4.3, 3.4.5, 3.4.7 and 3.4.10 we can show the existence of a minimal degree from the following:

Lemma 3.4.12 Let $T$ be a computable tree and let $e$ be an index. Then there is a computable subtree $S$ that is either $e$-collecting or $e$-splitting.

Proof
Case 1: There is a sequence $\sigma$ such that for all $\tau$ and $\tau'$ extending $\sigma$, $\phi^T_e(\tau)$ and $\phi^T_e(\tau')$ are equal where both are defined. Let $S(\tau) = T(\sigma * \tau)$. Then $S$ will be a subtree of $T$ and $S$ will be $e$-collecting.
Case 2: Otherwise. Then for every $\sigma$ there will be extensions $\tau_0$ and $\tau_1$ such that $\phi^T_e(\tau_0)$ and $\phi^T_e(\tau_1)$ are inconsistent. Further, we can find these $\tau_0$ and $\tau_1$ as computable functions $t_0(\sigma)$ and $t_1(\sigma)$.
We then define the subtree $S$ by

$$S(() ) = T(() ),$$

$i.e. S and T are equal on the empty sequence.$

If $S(\sigma)$ is defined, let $S(\sigma * 0) = T(t_0(\sigma))$ and $S(\sigma * 1) = T(t_1(\sigma))$.
This defines a subtree that will be $e$-splitting.

We have now proved all essential steps needed in the construction of a minimal degree:

Theorem 3.4.13 There is a non-computable function $f$ such that if $g <_T f$ then either $g$ is computable or $g$ is equivalent to $f$.

Proof
Using the lemmas above we construct a family $\{T_n\}$ of computable trees such that $T_{n+1}$ is a subtree of $T_n$ for all $n$, and such that for all $e$:

- If $\phi^e_e$ is total, then $\phi^e_e$ is not a branch in $T_{2e+1}$.
- $T_{2e+2}$ is either $e$-collecting or $e$-splitting.

Then by Lemma 3.4.3 there is an $f$ that is a branch in all $T_n$’s, and this $f$ will have the property wanted.
3.5 A priority argument

3.5.1 C.e. degrees

In the constructions of minimal pairs and functions of minimal degrees we have not been concerned with the complexity of the sets and functions constructed. We can decide upper bounds on the complexity of the objects constructed in the proofs by analyzing the complexity of properties like \( \phi_e \text{ is total} \) and counting in depth how many number quantifiers we will need in order to write out a definition of the object constructed. If we, however, are interested in results about degrees with some bounded complexity, we must be more careful in our constructions. In this section we will be interested in degrees with at least one c.e. set in it:

**Definition 3.5.1** Let \( a \) be a degree.

\( a \) is an c.e. degree if \( a \) contains the characteristic function of a c.e. set. We say that \( f \) is of c.e. degree if the degree of \( f \) is a c.e. degree.

There is a nice characterization of the functions of c.e. degree. We leave the proof as an exercise for the reader, see Exercise 3.26.

**Theorem 3.5.2** Let \( f \) be a function. Then the following are equivalent:

i) \( f \) is of c.e. degree.

ii) There is a primitive recursive sequence \( \{f_i\}_{i \in \mathbb{N}} \) converging pointwise to \( f \) such that the following function

\[ g(x) = \mu n. \forall m \geq n (f_m(x) = f(x)) \]

is computable in \( f \).

3.5.2 Post's Problem

So far we only know two c.e. degrees, \( O \), the degree of the computable sets and functions, and \( O' \), the degree of the halting problem or of \( K \). Post's Problem asks if there are more c.e. degrees than those two. This is of course a nice, technical problem, but it has implications beyond that. One of the reasons why c.e. sets are so interesting is that the set of theorems in an axiomatizable theory is c.e. If there were no more c.e. degrees than those two known to us, a consequence would be that there are two kinds of axiomatizable theories, those that are decidable and those that share the complexity of Peano Arithmetic.

As a consequence of Gödel's proof of the incompleteness theorem, the set of Gödel-numbers of theorems in Peano Arithmetic is a complete c.e. set, and actually of even the same \( m \)-degree as \( K \).

Now, in 1957 two young mathematicians, Friedberg and Muchnic, independently constructed c.e. sets of in-between degrees. They both developed what is now known as the priority method. The problem we have to face when constructing c.e. sets is that we must give an algorithm for adding elements to the...
set, but we cannot give an algorithm for keeping objects out of the set. If we did that, the set constructed would become computable. Thus when we have made an attempt to approximate a set with positive and negative information, we must be allowed to violate the negative information. However, we must not violate the negative information to such an extent that we ruin our global goal. We solve this dilemma by introducing priorities to our requirements. If an effort to satisfy one requirement will ruin the attempt to satisfy another requirement, we let the requirement with highest priority win. This idea will work when two properties are satisfied by the construction: If we make an attempt to satisfy a requirement and we never ruin this attempt, we actually manage to satisfy the requirement. Further, if we after a stage in the construction never make an attempt to satisfy a requirement, the requirement will automatically be satisfied.

Thus we are bound to satisfy the requirement of highest priority, either because we make an attempt which will not be ruined, or because there is no need to make an attempt. Then we are bound to satisfy the next requirement, either because we make an attempt after the final attempt for the first requirement, or because there is no need to make such an attempt, and so on... In the finite injury lemma we will give a full proof along this line of arguing.

### 3.5.3 Two incomparable c.e. degrees

Post’s problem was solved by constructing two incomparable c.e. degrees $a$ and $b$. We will see below why this actually solves the original problem.

**Theorem 3.5.3** There are two c.e. sets $A$ and $B$ that are not computable relative to each other.

**Remark 3.5.4** If $A$ and $B$ are not computable in each other, neither can be computable, because any computable set will be computable in any set. Moreover neither can have the same degree as $K$, because every c.e. set is computable in $K$. Thus we have not just produced an in-between degree, but two in-between degrees. In Exercise 3.27 we will see that there are infinitely many in-between c.e. degrees, and that any countable partial ordering can be embedded into the ordering of the c.e. degrees.

**Proof**

We will construct two c.e. sets $A$ and $B$ satisfying

- $R_{2e}: \ N \setminus A \neq W_e^B$
- $R_{2e+1}: \ N \setminus B \neq W_e^A$

or in other terms: The complement of $A$ is not c.e. relative to $B$ and vice versa. If we achieve this for all $e$, we will have proved the theorem, since a set $A$ is

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computable in $B$ if and only if both $A$ and the complement of $A$ are c.e. in $B$, and since $A$ is c.e. we have that $A$ is computable in $B$ if and only if the complement of $A$ is c.e. in $B$.

We will construct two primitive recursive increasing sequences $\{A_n\}_{n \in \mathbb{N}}$ and $\{B_n\}_{n \in \mathbb{N}}$ of finite sets. We let

$$A_0 = B_0 = \emptyset.$$ 

We call each step in the process a stage. If $n = \langle 1, e, x \rangle$ we will consider to make an attempt to satisfy $R_{2e}$ at stage $n$, while if $n = \langle 2, e, x \rangle$ we will consider to make an attempt to satisfy $R_{2e+1}$.

An attempt to satisfy $R_{2e}$ will consist of selecting a $q \in A_{n+1} \cap W_e^{B_{n+1}}$, and then put up a protection, the set of points used negatively in the verification of $q \in W_e^{B_{n+1}}$. If we can keep all objects in this protection out of $B$ throughout the construction, we will have

$$q \in A \cap W_e^B$$

and $R_{2e}$ will be satisfied. We call the protection active at a later stage $m$ if $B_m$ is disjoint from this protection.

There is some little minor trick to observe, we will use disjoint infinite supplies of numbers that we may put into $A$ (or $B$) in order to satisfy $R_{2e}$ (or $R_{2e+1}$). We will use this to show that if we make only finitely many attempts, we will succeed after all.

Now let $n = \langle 1, e, x \rangle$ and assume that $A_n$ and $B_n$ are constructed. Assume further that we constructed certain protections, some of them active at stage $n$, others not. We write the following procedure for what to do next:

Let $B_{n+1} = B_n$.

**Question 1:** Is there a protection for $R_{2e}$ active at stage $n$?

If the answer is ‘yes’, let $A_{n+1} = A_n$, and continue to the next stage.

If the answer is ‘no’, ask

**Question 2:** Is there a $y < n$ such that $\phi^{B_n}_{e,n}(\langle y, e \rangle) \downarrow$ and $y$ is in no active protection for any requirement $R_{2d+1}$ where $2d + 1 < 2e$?

If the answer is ‘no’, let $A_{n+1} = A_n$ and proceed to the next stage.

If the answer is ‘yes’, choose the least $y$, let $A_{n+1} = A_n \cup \{y, e\}$, construct a protection $\{0, \ldots, n\} \setminus B_n$ for $R_{2e}$ and move on to the next stage.

If $n = \langle 2, e, x \rangle$ we act in the symmetric way, while for other $n$ we just move on to the next stage, not adding anything to $A$ or $B$.

This ends the construction.

**Claim 1 (The Finite Injury Lemma)**

For each requirement $R_s$ there is a stage $n_s$ after which we do not put up or injure any protection for $R_s$.

**Proof**

We prove this by induction on $s$, and as an induction hypothesis we may assume that there is a stage $m_s$ after which we never put up a protection for any requirement $R_t$ with $t < s$. Then after stage $m_s$ we will never injure a protection.
for $R_s$. Thus if we never put up a protection for $R_s$ after stage $m_s$ we can let $n_s = m_s$, while if we construct a protection, this will never be injured and we can let $n_s$ be the stage where this protection is constructed.

Now let $A = \bigcup_{n \in \mathbb{N}} A_n$ and $B = \bigcup_{n \in \mathbb{N}} B_n$. Then $A$ and $B$ are c.e. sets.

Claim 2
Each requirement $R_s$ will be satisfied.

Proof
We prove this for $s = 2e$. There are two cases.
1. There is a protection for $R_s$ active at stage $n_s$.
   If this is the case, there will be a $y$ such that $\langle y, e \rangle \in A_{n_s} \cap W_e^{B_e}$. Since the protection is not injured, we will have that $\langle y, e \rangle \in A \cap W_e^B$ and the requirement is satisfied, $A$ and $W_e^B$ are not complementary.
2. There is no such protection.
   There are only finitely many objects of the form $\langle y, e \rangle$ in $A$, because we add at most one such object for each stage before $n_s$, and never any at a stage after $n_s$.

On the other hand, there can be only finitely many objects of the form $\langle y, s \rangle$ in $W_e^B$, since otherwise we could choose one that is not in any protection for any $R_t$ for $t < s$, and sooner or later we would at some stage after $n_s$ make a new attempt to satisfy $R_e$, which we are not. Thus $A \cup W_e^B$ contains only finitely many objects of the form $\langle y, e \rangle$ and the sets are not the complements of each other. Thus the requirement will be satisfied in this case as well.

We have shown that all the requirements are satisfied in this construction, so the theorem is proved.

3.6 Subrecursion theory

3.6.1 Complexity
What is to be considered as complex will be a matter of taste. Actually, the same logician may alter her/his taste for complexity several times a day. The generic logician may give a class on automata theory in the morning, and then the regular languages will be the simple ones, while context free languages are more complex. Still, context free languages are decidable in polynomial time, and while our logician spends an hour contemplating on the $P = NP$-problem any context free language is by far simpler than the satisfiability problem for propositional logic. If our logician is working mainly in classical computability theory, all decidable languages are simple, while the undecidable ones are the complex ones. If, however, our logician is a set theorist, all definable sets are simple, maybe all subsets of subsets of subsets of $\mathbb{R}$ are simple, we have to move high up in cardinality in order to find sets of a challenging and interesting complexity.

In this section, our view on complexity will be one shared by many proof theorists. One of the aims in proof theory is to characterize the functions and sets provably computable in certain formal theories extending elementary number
theory. The idea is that if we know the functions provably computable in \( T \), we know something worth knowing about the strength of the theory \( T \).

We will not be concerned with proof theory in this compendium, and any references to proof-theoretical results should not be considered as a part of any curriculum based on this text.

### 3.6.2 Ackermann revisited

In Section 3.2.1 we defined the Ackermann branches. The idea of Ackermann was that each use of the scheme for primitive recursion involves an *iteration* of a previously defined function. Then diagonalising over a sequence of functions defined by iterated iteration would break out of the class of primitive recursive functions.

The Ackermann branches were defined using functions of two variables. We will be interested in pushing his construction through to the transfinite level, in order to describe more complex functions from below. Then it will be convenient, from a notational point of view, to use functions of one variable.

**Definition 3.6.1**
- Let \( F_0(x) = x + 1 \)
- Let \( F_{k+1}(x) = F_{x+2}^{x+2}(x) \)

In Exercise 3.30 we will see that diagonalising over the \( F_k \)'s will lead us outside the class of primitive recursive functions.

From now on in this section we will let \( PA \) be first order Peano arithmetic; i.e. elementary number theory with the axiom scheme for first order induction. Our first order language \( L \) will be the language of \( PA \).

**Definition 3.6.2** A function \( f : \mathbb{N} \to \mathbb{N} \) is provably computable if there is a formula \( A(x,y) = \exists z B(x,y,z) \) in \( L \) where \( A \) defines the graph of \( f \), \( B \) has only bounded quantifiers and \( PA \vdash \forall x \exists y A(x,y) \).

This definition is extended in the obvious way to functions of several variables.

We will assume that the reader is familiar with the method of arithmetisation, sequence numbering and so forth.

**Lemma 3.6.3** Let \( f(k,x) = F_k(x) \). Then \( f \) is provably computable.

**Outline of proof**
Essentially we have to prove that \( F_0 \) is total and that if \( F_k \) is total then \( F_{k+1} \) is total as well. This induction step is proved by induction, where we actually must prove

\[
\forall x \forall y \exists z F_k^y(x) = z
\]

by induction on \( y \).

In order to write a complete proof, we have to write down the formal definition
of the graph of \( F \) and the establish a proof tree for the formula required to be a theorem in \( PA \). This is tedious and not very exiting; the exiting part is to decide how much induction that is required.

### 3.6.3 Ordinal notation

The concept of an ordinal number will be properly introduced in a course on set theory. We have defined the concept of well orderings. An ordinal number will be a set that in a canonical way represents an isomorphism class of well orderings. For our purposes it will be sufficient to think about order types of well orderings, and we will use a term language for such order types. Exercise 5.11 is a self-service introduction to ordinal numbers.

**Definition 3.6.4** \( \omega \) is the smallest infinite ordinal number, and it denotes the order-type of \( \mathbb{N} \) with the standard ordering.

All natural numbers will be ordinals numbers as well. We may extend the arithmetical operations *plus*, *times* and *exponentiation* to operations on orderings, or actually, on order types. The sum of two orderings \( \langle A, <_A \rangle \) and \( \langle B, <_B \rangle \) will be the set

\[
C = A \oplus B = (\{0\} \times A) \cup (\{1\} \times B)
\]

with the ordering \( <_C \) defined by

- \( \langle 0, a \rangle <_C \langle 1, b \rangle \) whenever \( a \in A \) and \( b \in B \).
- \( \langle 0, a_1 \rangle <_C \langle 0, a_2 \rangle \) if and only if \( a_1 <_A a_2 \).
- \( \langle 1, b_1 \rangle <_C \langle 1, b_2 \rangle \) if and only if \( b_1 <_B b_2 \).

The product of two orderings will for our purposes be the anti-lexicographical ordering on the product of the domains.

The formal definition of the exponential of orderings is less intuitive, but it helps to think of exponentiation as iterated multiplication. Our definition only works for special orderings \( \langle A, <_A \rangle \), including all well orderings.

**Definition 3.6.5** Let \( \langle A, <_A \rangle \) and \( \langle B, <_B \rangle \) be two orderings where \( A \) has a least element \( a_0 \).

Let \( p \in C \) if \( p \) is a map from \( B \) to \( A \) such that \( p(b) = a_0 \) for all but finitely many \( b \in B \).

If \( p \neq q \) are in \( C \), there will be a maximal argument \( b \) such that \( p(b) \neq q(b) \).

We then let \( p <_C q \Leftrightarrow p(b) < q(b) \) for this maximal \( b \).

**Lemma 3.6.6** If \( \langle A, <_A \rangle \) and \( \langle B, <_B \rangle \) are two well orderings then the sum, product and exponential will be well orderings.
The proof is left as Exercise 3.31

As a consequence, every arithmetical expression in the constant $\omega$ and constants for the natural numbers, and using 'plus', 'times' and 'exponents', will have an interpretation as an ordinal number. The least ordinal number that cannot be described by an expression as above is baptized $\epsilon_0$. As for ordinary arithmetics, $\omega^0 = 1$. This is a special case of the more general rule

$$\omega^\alpha \cdot \omega^\beta = \omega^{\alpha+\beta}.$$

A consequence is that each ordinal $\alpha < \epsilon_0$ can be written in a unique way as

$$\alpha = \omega^{\alpha_n} + \cdots + \omega^{\alpha_0}$$

where $\{\alpha_0, \ldots, \alpha_n\}$ is an increasing (not necessarily strictly) sequence of ordinals less than $\alpha$. We call this the Cantor Normal Form of $\alpha$. If $\alpha_0 = 0$ the ordinal $\alpha$ will be a successor ordinal, otherwise it will be a limit ordinal.

We extended the Ackermann hierarchy to the first transfinite level by diagonalising over the Ackermann branches. One advantage with the Cantor normal form is that we may find a canonical increasing unbounded sequence below every limit ordinal between $\omega$ and $\epsilon_0$. Actually, it is possible to do so for ordinals greater than $\epsilon_0$ too, but readers interested in how this is done and why someone would like to do it are recommended to follow a special course on proof theory and ordinal denotations.

**Definition 3.6.7** Let

$$\alpha = \omega^{\alpha_n} + \cdots + \omega^{\alpha_0}$$

be given on Cantor normal form, and assume that $\alpha_0 > 0$.

We define the $n$'th element $\alpha[n]$ of the fundamental sequence for $\alpha$ as follows:

**Case 1** $\alpha_0 = \beta + 1$:

Let $\alpha[n] = \omega^{\alpha_m} + \cdots + \omega^{\alpha_1} + n \cdot \omega^\beta$.

**Case 2** $\alpha_0$ is a limit ordinal:

We may then assume that $\alpha_0[n]$ is defined, and we let

$$\alpha[n] = \omega^{\alpha_m} + \cdots + \omega^{\alpha_1} + \omega^{\alpha_0}[n]$$

We may consider the map $\alpha \mapsto \alpha[n]$ as an extension of the predecessor function to limit ordinals:

**Definition 3.6.8** Let $0 < \alpha < \epsilon_0$. We define the $n$-predecessor of $\alpha$ by

a) If $\alpha$ is a successor ordinal, the predecessor of $\alpha$ will be the $n$-predecessor of $\alpha$.

b) If $\alpha$ is a limit ordinal, $\alpha[n]$ will be the $n$-predecessor of $\alpha$. 
c) We say that $\beta <_n \alpha$ if $\beta$ can be reached from $\alpha$ by iterating the $n$-predecessor map.

**Lemma 3.6.9** Let $\alpha < \epsilon_0$, $\beta <_m \alpha$ and $m < n$. Then $\beta <_n \alpha$.

**Proof**
It is sufficient to prove that if $m < n$ and $\alpha$ is a limit ordinal, then $\alpha[m] <_n \alpha$.
This is left as a non-trivial exercise for the reader, see Exercise 5.7.

**Lemma 3.6.10** Let $\alpha < \epsilon_0$ and let $\beta < \alpha$. Then there is an $n$ such that $\beta <_n \alpha$.

**Proof**
This is proved by induction on $\alpha$ with the aid of Lemma 3.6.9. The details are left as a nontrivial exercise for the reader, see Exercise 5.7.

### 3.6.4 A subrecursive hierarchy

We will now extend the alternative Ackermann hierarchy to all ordinals less than $\epsilon_0$.

**Definition 3.6.11** Let $\alpha < \epsilon_0$. We define $F_\alpha(x)$ by recursion on $\alpha$ as follows:
- $F_0(x) = x + 1$
- $F_{\beta+1}(x) = F_{\alpha+2}(x)$
- $F_\alpha(x) = F_\alpha[x](x)$ when $\alpha$ is a limit ordinal.

Proof theorists have shown that if a function $f : \mathbb{N} \to \mathbb{N}$ is provably computable, then $f$ will be bounded almost everywhere by one of the $F_\alpha$’s where $\alpha < \epsilon_0$.

The proof involves the translation of a proof in $PA$ to a proof in $\omega$-logic, then a cut-elimination procedure in $\omega$-logic and finally an analysis of cut-free proofs in $\omega$-logic of the totality of computable functions. This analysis is far beyond the scope of this compendium, and the reader is not recommended to work out the details.

In Exercises 5.7 and 5.8, the reader is challenged to prove that each $F_\alpha$ is provably computable and establish the hierarchical properties of the $\{F_\alpha\}_{\alpha < \epsilon_0}$ hierarchy. The main results are:

**Lemma 3.6.12** Let $n < x$ and $\beta <_n \alpha < \epsilon_0$.

Then

$$F_\beta(x) \leq F_\alpha(x).$$

**Theorem 3.6.13** If $\beta < \alpha < \epsilon_0$, then

$$\exists x_0 \forall x > x_0 (F_\beta(x) < F_\alpha(x)).$$

**Remark 3.6.14** There are many aspects about subrecursive hierarchies that we have not discussed in this section. We have not discussed complexity classes.

For instance, the class $H_\alpha$ of functions computable in polynomial time relative to a finite iteration of $F_\alpha$ represents a stratification of the set of all provably computable functions into a complexity hierarchy, where each complexity class is closed under composition and closed under polynomial time reductions. We will not discuss such matters further here.

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3.7 Exercises

Exercise 3.1 Prove that the following functions are primitive recursive:

a) \( f(x, y, \vec{z}) = x + y \)

b) \( f(x, y, \vec{z}) = x \cdot y \)

c) \( f(x, y, \vec{z}) = x^y \)

d) \( f(x, \vec{y}) = x! \)

e) \( f(x, \vec{y}) = x - 1 \)

f) \( f(x, y, \vec{z}) = x - y \)

g) \( f(x, y, \vec{z}) = 0 \) if \( x = y \)
\( f(x, y, \vec{z}) = 1 \) if \( x \neq y \)

h) \( f(x, y, \vec{z}) = 0 \) if \( x < y \)
\( f(x, y, \vec{z}) = 1 \) if \( x \geq y \)

Exercise 3.2 Prove the following facts:

a) \( \mathbb{N}^n \) and \( \emptyset \) are primitive recursive as subsets of \( \mathbb{N}^n \).

b) The complement of a primitive recursive set is primitive recursive. Moreover, the union and intersection of two primitive recursive subsets of \( \mathbb{N}^n \) will be primitive recursive.

Exercise 3.3

a) Prove Lemma 3.2.6.

b) Prove that we can replace the inequalities by strict inequalities in Lemma 3.2.6.

Exercise 3.4 Prove Lemma 3.2.7.

Exercise 3.5

a) Prove Lemma 3.2.11.

b) Prove that the sequence numbering is monotone in each coordinate.

c) Prove that the monotone enumeration \( SEQ \) of the sequence numbers is primitive recursive.

Hint: Find a primitive recursive bound for the next sequence number and use bounded search.

d) Define an alternative sequence numbering as follows:
\( \langle x_0, \ldots, x_{n-1} \rangle \) is the number \( z \) such that
\[
SEQ(z) = \langle x_0, \ldots, x_{n-1} \rangle.
\]
Show that this alternative numbering is surjective and still satisfies Lemma 3.2.11.
e) Prove that the pairing function $P$ in Definition 3.2.13 is 1-1 and onto.

**Exercise 3.6** Let $X$ be a set of functions closed under the schemes of primitive recursion.
Show that for any function $f : \mathbb{N} \rightarrow \mathbb{N}$ we have

$$f \in X \iff \bar{f} \in X.$$  

**Exercise 3.7** We define the function $f(k, e, y)$ by recursion on $k$ and subrecursion on $e$ as follows:

1. For all $k$, if $e = (1)$ and $y = \langle x, x_1, \ldots, x_n \rangle$ we let
   $$f(k, e, y) = x + 1.$$  

2. For all $k$, if $e = (2, i)$, $y = \langle x_1, \ldots, x_n \rangle$ and $1 \leq i \leq n$ we let
   $$f(k, e, y) = x_i.$$  

3. For all $k$, if $e = (3, q)$ we let
   $$f(k, e, y) = q.$$  

4. For all $k$, if $e = (4, e', d_1, \ldots, d_n)$ we let
   $$f(k, e, y) = f(k, e', \langle f(k, d_1, y), \ldots, f(k, d_n, y) \rangle).$$  

5. For all $k > 0$, if $e = \langle 5, e_1, e_2 \rangle$ we let
   
   * $f(k, e, \langle 0, x_1, \ldots, x_n \rangle) = f(k - 1, e_1, \langle x_1, \ldots, x_n \rangle)$
   
   ** $f(k, e, \langle m + 1, x_1, \ldots, x_n \rangle) = f(k - 1, e_2, \langle f(k, e, \langle m, x_0, \ldots, x_n \rangle), m, x_1, \ldots, x_n \rangle)$
   
6. In all other cases, we let $f(k, e, y) = 0$.

a) Prove that $f$ is well defined, and that $f$ is computable.

b) Prove that if $g$ of arity $n$ is primitive recursive, there is a number $k$ and an index $e$ such that

$$g(x_1, \ldots, x_n) = f(k, e, \langle x_1, \ldots, x_n \rangle)$$

for all $x_1, \ldots, x_n \in \mathbb{N}^n$.

c) Prove that $f$ is not primitive recursive.

**Exercise 3.8** Prove that if $\phi_e(\bar{x}) = y$ and $\phi_e(\bar{x}) = z$, then $y = z$.

**Hint:** Use induction on $e$.

Discuss why this is something that needs a proof.
Exercise 3.9 Let \( \{f_n\}_{n \in \mathbb{N}} \) be a sequence of total functions such that
\[
g(n, m) = f_n(m)
\]
is computable.
Show that each \( f_n \) is computable, and that there is a total computable function not in the sequence.
Hint: Use a diagonal argument.

Exercise 3.10 Show that there is a total computable function \( \Phi(n, x) \) of two variables that enumerates all primitive recursive functions of one variable.
Is it possible to let \( \Phi \) be primitive recursive?

Exercise 3.11 Complete the proof of Lemma 3.2.27.

Exercise 3.12 Prove Corollary 3.2.32.

Exercise 3.13 a) Prove that every non-empty c.e. set is the image of a primitive recursive function (easy) and that every infinite c.e. set is the image of an injective computable function (not that easy, but still...).

b) Prove that the range of a strictly increasing total computable function is computable.

Exercise 3.14 Let \( A \subseteq \mathbb{N}^n \). Show that \( A \) is c.e. (by characterization ii) or iii) in Theorem 3.2.35) if and only if
\[
\{ \langle x_1, \ldots, x_n \rangle ; (x_1, \ldots, x_n) \in A \}
\]
is c.e.

Exercise 3.15 Give an explicit description of the selection function in the proof of Theorem 3.2.39.

Exercise 3.16 Let \( A \) and \( B \) be two disjoint sets whose complements are c.e. Show that there is a computable set \( C \) such that \( A \subseteq C \) and \( B \cap C = \emptyset \).
Hint: Use Corollary 3.2.40.

Exercise 3.17 Let \( L \) be the language of propositional calculus over an infinite set \( \{A_i\}_{i \in \mathbb{N}} \) of propositional variables. Discuss the following statement: There is a primitive recursive consistent set of propositions with no computable completion.

Exercise 3.18 a) Show that there is an enumeration \( \{I_n\}_{n \in \mathbb{N}} \) of all closed rational intervals contained in \([0, 1]\) such that the relations \( I_n \subseteq I_m \), \( I_n \cap I_m = \emptyset \) and \( |I_n| < 2^{-m} \) are computable, where \( |I_n| \) is the length of the interval.

A real number \( x \) is \textit{computable} if there is a computable function \( h \) such that
b) Let $f : [0, 1] \to [0, 1]$ be a continuous function. We say that $f$ is computable if there is a total computable function $g : \mathbb{N} \to \mathbb{N}$ such that

i) $I_n \subseteq I_m \Rightarrow I_{g(n)} \subseteq I_{g(m)}$

ii) $x \in I_n \Rightarrow f(x) \in I_{g(n)}$

iii) For all $x$ and $m$ there is an $n$ with $x \in I_n$ and $|I_{g(n)}| < 2^{-m}$.

Show that any computable function will be continuous. Show that there is a computable function that does not take its maximal value at any computable real.

Exercise 3.19 Show that if $A \equiv_m B$ and $C \equiv_m D$, then $A \oplus C \equiv_m B \oplus D$. Show that if $A <_m E$ and $B <_m E$ then $A \oplus B <_m E$. Show that $A \oplus B$ then represents the least upper bound of $A$ and $B$ in the $m$-degrees.

Exercise 3.20 Show that the second recursion theorem can be extracted from the proof of the recursion theorem.

Hint: Let $F(g, x)$ be computable. Use the recursion theorem on $f(e, x) = F(\phi_e, x)$.

Let $e_0$ be the index obtained from the proof of the recursion theorem. Show that if $\phi_{e_0}(x) = y$ and we make an oracle call for $\phi_{e_0}(z)$ in the computation of $F(\phi_{e_0}, x)$ then $\phi_{e_0}(z)$ is a subcomputation of $\phi_{e_0}(x)$.

Exercise 3.21 Show that if $f_1 \equiv_T f_2$, then the jumps $f_1'$ and $f_2'$ are also Turing equivalent.

Exercise 3.22 Make a direct construction and show that there are two total functions $f$ and $g$ such that $f \not\equiv_T g$ and $g \not\equiv_T f$.

Exercise 3.23 Prove that there is a minimal pair of Turing degrees, i.e. a pair \{a, b\} of degrees of non-computable functions, such that the degree $O$ of computable functions is the greatest lower bound of $a$ and $b$.

Hint: You may use the main idea in the proof of Theorem 3.3.27, the present theorem is just simpler to prove. You may also get some ideas from the discussion of this proof in the text.

Exercise 3.24 Let $O$ be the degree of the computable functions. Recall the definition of the jump operator in the proof of Lemma 3.3.26, cfr. Exercise 3.21. We define the arithmetical hierarchy as follows:

A $\Sigma^0_1$-set is a c.e. set (of any dimension).
A $\Pi_k^0$-set is the complement of a $\Sigma_k^0$-set
A $\Sigma_k^{0,k+1}$-set is the projection of a $\Pi_k^0$-set
A $\Pi_k^{0,k+1}$-set is the complement of a $\Sigma_k^0$-set
A $\Delta_k^0$ set is a set that is both $\Sigma_k^0$ and $\Pi_k^0$.

Let $O^{(n)}$ be the degree obtained from $O$ using the jump-operator $n$ times.

a) Prove that if $A$ is $\Sigma_k^0$ or $A$ is $\Pi_k^0$, then (the characteristic function of) $A$ has a degree $a \leq O^{(k)}$.

b) Show that for $k \geq 1$ we have that $A$ is $\Delta_k^0$ if and only if the degree of $A$ is bounded by $O^{(k)}$.

Hint: Use the relativized version of the fact that a set is computable if and only if both the set and its complement are c.e.

Exercise 3.25 Show that there are continuumly many minimal degrees.
Hint: Instead of constructing one tree $T_n$ at level $n$ we might construct $2^n$ trees $T_{\sigma,n}$ for $h(\sigma) = n$, ensuring for each $e$ that the branches of different trees will not be computable in each other via index $e$. The proof requires some clever book-keeping.

Exercise 3.26 Prove Theorem 3.5.2.

Exercise 3.27 Let $B \subseteq \mathbb{N}^2$ be a set. For each $n \in \mathbb{N}$, we let $B_n = \{m : (n,m) \in B\}$ and we let $B_{-n} = \{(k,m) \in B : k \neq n\}$.

a) Show that there is a c.e. set $B$ such that for all $n$, $B_n$ is not computable in $B_{-n}$.

Hint: Use an enumeration of $\mathbb{N}^2$ to give all the requirements

$$R_{(n,c)} : \mathbb{N} \setminus B_n \neq W_{c,b}^{B_{-n}}$$

a priority rank.

b) Consider a computable partial ordering $\prec$ on the natural numbers.

Show that there is an order-preserving map of $\prec$ into the c.e. degrees.

Hint: Use the construction in a), and let $C_n = \{(k,m) \in B : k \preceq n\}$.

There is one computable partial ordering $\prec$ such that any other partial ordering of any countable set can be embedded into $\prec$, see Exercise 5.2. Thus this shows that any countable partial ordering can be embedded into the ordering of the c.e. degrees.

Exercise 3.28 Fill in the details in the proof of the following theorem:

Theorem 3.7.1 Let $a > O$ be an c.e. degree. Then there are two incomparable c.e. degrees $b$ and $c$ such that $a = b \oplus c$. 

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This theorem is called \textit{The splitting theorem}. We split the c.e. degree \(a\) into two simpler c.e. degrees.

\textbf{Proof}

Let \(A\) be a non-computable c.e. set. It is sufficient to construct two disjoint c.e. sets \(B\) and \(C\) such that \(A = B \cup C\), \(A\) is not computable in \(B\) and \(A\) is not computable in \(C\).

Let \(f\) be a 1-1 enumeration of \(A\). At each stage \(n\) we will put \(f(n)\) into \(B\) or \(f(n)\) into \(C\), but not into both.

Let \(B_n\) be the numbers put into \(B\) before stage \(n\) and \(C_n\) be the set of numbers put into \(C\) before stage \(n\). Further, we let \(A_n = \{f(0), \ldots, f(n-1)\}\), so \(A_n = B_n \cup C_n\).

We put up requirements \(R_{2e}\) :

\[ K_A \neq \phi^B_e. \]

\(R_{2e+1} : K_A \neq \phi^C_e. \)

which we give priorities in the usual way.

For each requirement \(R_s\) we define three auxiliary functions. For \(s = 2e\) they will be:

\textit{The match function}

\[ m(s, n) = \mu k < n. \forall x \leq k(\phi^B_{e,n}(x) = K_{A_n}(x)). \]

\textit{The bar function}

\[ b(s, n) = \max \{m(s, n') : n' \leq n\}. \]

\textit{The protection function}

\[ p(s, n) = \{y : y \text{ is used negatively in computing } \phi^B_{e,n}(x) \text{ for some } x \leq b(s, n)\}. \]

In this case we call this \textit{a protection of }\(B\).

Now the construction at stage \(n\) is as follows: If \(f(n) \not\in p(s, n)\) for any \(s \leq n\), put \(f(n)\) into \(B\). Otherwise, consider the requirement \(R_s\) of highest priority such that \(f(n) \in p(s, n)\). If this is a protection of \(B\), we put \(f(n)\) into \(C\), otherwise we put \(f(n)\) into \(B\).

When we put an element into a protection of \(B\) we \textit{injure} that requirement. We will prove that for any requirement \(R_s\) there is a stage \(n_s\) after which we will never injure that requirement, and simultaneously that the bar-function \(b(s, n)\) is bounded when \(s\) is fixed.

Assume that this holds for all \(s' < s\). Then there is a stage \(n_s\) after which \(f(n)\) is not in the protection for any \(s' < s\), and then, after stage \(n_s\), \(R_s\) will not be injured.

This in turn means that if \(x < b(s, n)\) for \(n \geq n_s\) and \(\phi^B_{e,n}(x)\), then \(\phi^B(x) = \phi^B_{e,n}(x)\).

Now, if \(\lim_{n \to \infty} v(s, n) = \infty\) we can use the increasing matching and the stability
of \( \phi_{c,n} \) to show that \( A \) is computable, which it is not.
On the other hand, if \( K_A = \phi_c^B \) we will get increasing matching. Thus at
the same time we prove that the construction of the bar and protection for \( R_s \)
terminates and that the requirement is satisfied at the end.

**Exercise 3.29** Post hoped to prove that a simple set cannot be of the same
degree as the complete c.e. set \( K \). This will not be the case, which will be clear
when you have solved this problem.
Let \( A \) be a non-computable c.e. set and \( f \) a total computable 1-1-enumeration
of \( A \). We know that \( f \) cannot be increasing (why?).
Let
\[
B = \{ n : \exists m > n (f(m) < f(n)) \}
\]
This is called the *deficiency set* of the enumeration.

a) Show that \( B \) is c.e. and that \( B \) is computable n \( A \).
b) Show that \( A \) is computable in \( B \).
   Hint: In order to determine if \( x \in A \) it is sufficient to find \( n \not\in B \) such
   that \( f(n) > x \).
c) Show that \( B \) is simple.
   Hint: If the complement of \( B \) contains an infinite c.e. set, the algorithm for
   computing \( A \) from \( B \) in b) can be turned into an algorithm for computing
   \( A \).

**Exercise 3.30** Let \( F(k, x) = F_k(x) \) be the alternative Ackermann function.
a) Show that \( x < F(k, x) \) for all \( k \) and \( x \).
b) Show that \( F \) is monotone in both variables.
c) Let \( f : \mathbb{N}^m \to \mathbb{N} \) be primitive recursive.
   For \( \vec{x} = (x_1, \ldots, x_m) \), let \( \sum \vec{x} = x_1 + \cdots + x_m \). Show that there is a
   number \( k \) such that
   \[
f(\vec{x}) \leq F(k, \sum \vec{x})
\]
   for all \( \vec{x} \in \mathbb{N}^m \).
   Hint: Use induction on the construction of \( f \).
   In the cases \( f(x, \vec{x}) = x + 1 \) and \( f(\vec{x}) = x \), you may use \( k = 0 \).
   In the case \( f(\vec{x}) = q \), use a \( k \) such that \( F(k, 0) \geq q \).
   In the case of composition, you may find it convenient to increase \( k \) by
   more than one, while in the case of primitive recursion, you should increase
   \( k \) by exactly one.
d) Show that \( F \) is not primitive recursive.
   Hint: Show that \( G(k) = F(k, k) + 1 \) cannot be primitive recursive, using
c).

**Exercise 3.31** Prove Lemma 3.6.6
Exercise 3.32 Let $\mathcal{P}$ be the set of polynomials $P(x)$ where we only use $+$ (not 'minus') and where all coefficients are natural numbers. We order $\mathcal{P}$ by

$$P(x) \prec Q(x) \iff \exists n \forall m \geq n (P(m) < Q(m)).$$

Show that $\prec$ is a well ordering of order-type $\omega^\omega$. 
Chapter 4

Generalized Computability Theory

4.1 Computing with function arguments

When we introduced the Turing degrees, we first introduced the notion of relativized computations via the notation
\[ \phi_{e}^{f_1, \ldots, f_n}(x_1, \ldots, x_m). \]

By a small change of notation, we may view this as a partial functional
\[ \phi_{e}(x_1, \ldots, x_m, f_1, \ldots, f_n) \]
where some of the inputs may be numbers and some may be functions.

**Definition 4.1.1** Let \( F : \mathbb{N}^m \times (\mathbb{N}^\mathbb{N})^n \rightarrow \mathbb{N} \).

We say that \( F \) is computable if there is an index \( e \) such that for all \( \vec{f} \in (\mathbb{N}^\mathbb{N})^n \) and \( \vec{f} \in \mathbb{N}^m \) we have that
\[ F(\vec{x}, \vec{f}) = \phi_{e}(\vec{x}, \vec{f}). \]

We will concentrate our attention to computable functionals \( F : \mathbb{N}^\mathbb{N} \rightarrow \mathbb{N} \). Let us first see why we may expect to face all relevant theoretical problems even then.

**Lemma 4.1.2** Let \( n > 0 \) and \( m \geq 0 \). Then there is a computable bijection between \( \mathbb{N}^m \times (\mathbb{N}^\mathbb{N})^n \) and \( \mathbb{N}^\mathbb{N} \).

**Proof**

We will produce the bijections between \( (\mathbb{N}^\mathbb{N})^2 \) and \( \mathbb{N}^\mathbb{N} \) and between \( \mathbb{N} \times \mathbb{N}^\mathbb{N} \) and \( \mathbb{N}^\mathbb{N} \). The rest then follows by iterating the two constructions.

Let \( \langle f, g \rangle(x) = (f(x), g(x)) \) where we use a computable pairing function with
computable projections on \( \mathbb{N} \). Clearly \( \langle f, g \rangle \) is computable from \( f \) and \( g \) in the sense that
\[
F(x, f, g) = (f, g)(x)
\]
is computable. In the same sense, \( f \) and \( g \) will be computable from \( \langle f, g \rangle \).

If \( f \in \mathbb{N}^\mathbb{N} \) and \( x \in \mathbb{N} \), let
\[
\langle x, f \rangle(0) = x
\]
\[
\langle x, f \rangle(y + 1) = f(y)
\]
This clearly defines a bijection, and it is computable in the sense above.

4.1.1 Topology

All these spaces are actually topological spaces, even metrizable spaces. The topology on \( \mathbb{N} \) will be the discrete topology, and the natural metric will be the one where distinct numbers have distance 1.

The topology on \( \mathbb{N}^\mathbb{N} \) will be the standard product topology. Since there may be readers not familiar with the product topology, we give a direct definition.

Definition 4.1.3 Let \( O \subset \mathbb{N}^\mathbb{N} \). \( O \) will be open if whenever \( f \in O \) there is a number \( n \) such that for all \( g : \mathbb{N} \to \mathbb{N} \) we have that \( g \in O \) whenever \( f \) and \( g \) agree for all arguments between and including 0 and \( n - 1 \).

In more technical terms, this can be expressed by
\[
f \in O \Rightarrow \exists n \forall g(f(n) = g(n) \Rightarrow g \in O).
\]

If \( \sigma \) is a finite sequence of numbers, \( \sigma \) determines an open neighborhood
\[
B_\sigma = \{ f \in \mathbb{N}^\mathbb{N} ; \bar{f}(lh(\sigma)) = \sigma \}
\]
i.e. the set of functions \( f \) extending \( \sigma \). These sets \( B_\sigma \) will form a basis for the topology.

There is a close connection between topology and computability. There are many examples of this in the literature. Although the use of topology is restricted to fairly elementary general topology, it is very hard to read the current literature on computability on non-discrete structures without any knowledge of topology at all. In this section we will give one example of this connection.

Theorem 4.1.4 Let \( F : \mathbb{N}^\mathbb{N} \to \mathbb{N} \). Then the following are equivalent:

a) \( F \) is continuous.

b) There is a function \( f : \mathbb{N} \to \mathbb{N} \) such that \( F \) is computable relative to \( f \).
Proof
First let $F$ be computable relative to $f$, i.e. there is an index $e$ such that

$$F(g) = \phi_e(f,g)$$

for all $g$.

For a fixed $g$, the computation tree of $\phi_e(f,g)$ is finite, and thus application of $g$ will only be used a finite number of times in the computation tree. Let $n$ be so large that if $g(x)$ occurs in the computation tree, then $x < n$. As a consequence we see that if $h(n) = \tilde{g}(n)$ then the computation trees of $\phi_e(f,g)$ and $\phi_e(f,h)$ will be identical. The further consequence is that

$$\forall h \in \mathbb{N}(\tilde{h}(n) = \tilde{g}(n) \Rightarrow F(g) = F(h))$$

and this just means that $F$ is continuous.

Now, let us prove the converse, and let $F$ be continuous. Let $\{\sigma_n\}_{n \in \mathbb{N}}$ be a computable enumeration of all finite sequences. Let $X = \{n : F$ is constant on $B_{\sigma_n}\}$. Let $f$ be defined by

$$f(n) = 0 \text{ if } n \notin X$$

$$f(n) = m + 1 \text{ if } F \text{ is constant } m \text{ on } B_{\sigma_n}.$$ 

Then for each $g \in \mathbb{N} \rightarrow \mathbb{N}$ we have that

$$F(g) = f(\mu n. f(n) > 0 \land \sigma_n \prec g) - 1 \tag{4.1}$$

where $\sigma \prec g$ means that the finite sequence $\sigma$ is an initial segment of the infinite sequence $g$.

4.1.2 Associates

In the proof of Theorem 4.1.4 we constructed a function $f$ such that $f(n) > 0$ if and only if $F$ is constant on $B_{\sigma_n}$ and in this case, $f(n) = F(g) + 1$ for all $g \in B_{\sigma_n}$. Following Kleene, we call this $f$ the principal associate of $F$. This is of course an important concept, but from a computability theory point of view it is not completely satisfactory:

**Lemma 4.1.5** There is a computable functional $F$ such that the principal associate is not computable.

**Proof**

Let $A = W_e$ be c.e. but not computable, i.e.

$$n \in A \Leftrightarrow \exists m T(e, n, m)$$

where $T$ is Kleene’s $T$-predicate.

Define the computable $F$ by

- $F(g) = 0$ if $\exists m < g(1) T(e, g(0), m)$

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Let $f$ be the principal associate of $F$. Let $\nu(n) = \langle n \rangle$, i.e. the sequence number of the one-point sequence containing just $n$. $F$ will be constant on $B_{\langle n \rangle}$ if and only if $n \notin A$, and then the constant value is 1, so we have

$$n \notin A \Leftrightarrow f(\nu(n)) = 2.$$ 

Since $\nu$ is computable and $A$ is not computable, $f$ cannot be computable.

When we showed that any functional $F$ computable in $f$ will be continuous, we referred to the computation tree of $\phi_e(f,g)$. When we look at the example showing Lemma 4.1.5 we see that the computation tree for $F(g)$ always will make use of $g(0)$ and of $g(1)$. Let $f$ be defined by

- If $lh(\sigma_n) < 2$, let $f(n) = 0$.
- If $lh(\sigma_n) \geq 2$ and $\exists m < \sigma_n(1)T(e,\sigma_n(0),m)$, let $f(n) = 1$.
- If $lh(\sigma_n) \geq 2$ and $\forall m < \sigma_n(1)\neg T(e,\sigma_n(0),m)$, let $f(n) = 2$.

Then Equation 4.1 will hold for this $F$ and $f$.

This leads us to the following definition:

**Definition 4.1.6** Let $F : \mathbb{N}^\mathbb{N} \to \mathbb{N}$ be continuous. An associate for $F$ will be a function $f : \mathbb{N} \to \mathbb{N}$ such that

i) If $\sigma_n \prec \sigma_m$ and $f(n) > 0$ then $f(m) = f(n)$.

ii) If $f(\overline{g}(m)) > 0$ then $f(\overline{g}(m)) = F(g) + 1$.

iii) $\forall g \in \mathbb{N}^\mathbb{N}\exists m(f(\overline{g}(m)) > 0)$.

It is then clear that any continuous functional will be computable in any of its associates via equation 4.1, and any computable functional will have a computable associate, see Exercise 4.2

### 4.1.3 Uniform continuity and the Fan Functional

It is well known that any continuous function on a compact metric space is uniformly continuous. In $\mathbb{R}$, a set will be compact if and only if it is closed and bounded. We have a similar characterization for $\mathbb{N}^\mathbb{N}$:

**Definition 4.1.7**

a) We will consider the following partial ordering of $\mathbb{N}^\mathbb{N}$:

$$f \leq g \Leftrightarrow \forall n(f(n) \leq g(n))$$

b) For $f \in \mathbb{N}^\mathbb{N}$, let $C_f = \{ g : g \leq f \}$.

The following is left for the reader as Exercise 4.3
Lemma 4.1.8 Let $A \subseteq \mathbb{N}$. Then $A$ will be compact if and only if $A$ is closed and $A \subseteq C_f$ for some $f$.

As a result, each continuous $F$ will be uniformly continuous on each $C_f$. Let us see what this actually means. The formal definition is that for any $\epsilon > 0$ there is a $\delta > 0$ such that for all $g$ and $h$ in $C_f$, if $d(g, h) < \delta$ then $d(F(g), F(h)) < \epsilon$. The metric on $\mathbb{N}$ is trivial, so let $\epsilon = \frac{1}{2}$ and choose $\delta > 0$ accordingly. Choose $n$ such that $2^{n-1} < \delta$. Then $\bar{g}(n) = \bar{h}(n) \Rightarrow F(g) = F(h)$ whenever $g$ and $h$ are in $C_f$. These considerations give us

Lemma 4.1.9 Let $F : \mathbb{N}^N \rightarrow \mathbb{N}$ be continuous. Then

$$\forall f \in \mathbb{N}^F \exists n \forall g \in C_f \forall h \in C_f (\bar{g}(n) = \bar{h}(n) \rightarrow F(g) = F(h))$$

We suggest an alternative proof in Exercise 4.4.

Lemma 4.1.9 suggests that we may consider the operator $\Phi$ defined by:

Definition 4.1.10 The Fan Functional $\Phi$ is defined on all continuous $F : \mathbb{N}^N \rightarrow \mathbb{N}$ and defines a function $\Phi(F) : \mathbb{N}^N \rightarrow \mathbb{N}$ by

$$\Phi(F)(f) = \mu n. \forall g \in C_f \forall h \in C_f (\bar{g}(n) = \bar{h}(n) \rightarrow F(g) = F(h))$$

Theorem 4.1.11 Let $\Phi$ be the fan functional. Then $\Phi(F)$ is continuous whenever $F$ is continuous.

Proof

Let $f$ be given, and let $\alpha$ be an associate for $F$. For each $n$ there will be a finite set of sequences $\sigma$ of length $n$ that are bounded by $f$. Using König’s lemma we may find an $n$ such that $\alpha(\sigma) > 0$ for each $\sigma$ in this set. This $n$ can be found effectively in $f$ and $\alpha$, and will be an upper bound for $\Phi(F, f)$. We only need information from $\bar{f}(n)$ in order to verify that this $n$ is a good one, and we may compute $\Phi(F)(f)$ from the information at hand. The details are left as an exercise for the reader.

Remark 4.1.12 In some sense originally made precise by Kleene and independently by Kreisel, the fan functional is continuous, see Exercise 4.5. In the next section we will consider a concept of computation where we may accept even functionals as inputs. This begs the question if the fan functional is even computable. The answer to this is not unique, since there is no canonical choice of a concept of computability in this case. En Exercise 4.5 we give a positive answer to this question for one possible notion of computability.

We characterized the continuous functionals as those having associates. In a way this is an old fashioned approach. In the current literature one often use domains as an entry to the theory of continuous functionals. We did not do so, because then we would have to introduce a lot of simple, but new, concepts. A reader interested in learning more about this kind of generalized computability is advised to consult some introduction to domain theory.
4.2 Computing relative to a functional of type 2

A major step in the process of generalizing computability came when Kleene and others started to relativize computations to functionals of type 2 and higher types in general. For certain indices $e$, we have defined the partial function $\phi_e(x_1, \ldots, x_n, f_1, \ldots, f_m)$, and of course there is no harm in accepting functionals $F_1, \ldots, F_k$ as dummy arguments. The problem is how we, in any sensible way, can use $F_j$ actively in a computation.

The strategy will be that whenever we can supply $F_j$ with an argument $f$, then we get the value $F_j(f)$ from some oracle call scheme.

**Definition 4.2.1** Let $\Psi : \mathbb{N}^{n+1} \times (\mathbb{N}^n)^m \times (\mathbb{N}^n \to \mathbb{N})^k \to \mathbb{N}$ be given. Let $\vec{a} \in \mathbb{N}^n$, $\vec{f} \in (\mathbb{N}^n)^m$ and let $\vec{F} \in (\mathbb{N}^n \to \mathbb{N})^k$. By

$$\lambda x. \Psi(x, \vec{a}, \vec{f}, \vec{F})$$

we mean the function $g : \mathbb{N} \to \mathbb{N}$ defined by

$$g(x) = \Psi(x, \vec{a}, \vec{f}, \vec{F}).$$

This makes sense even when $\psi$ is a partial functional, but then $\lambda x. \Psi(x, \vec{a}, \vec{f}, \vec{F})$ may not be total.

Kleene suggested something equivalent to

**Definition 4.2.2** We extend the definition of $\phi_e(\vec{a}, \vec{f})$ to a definition of $\phi_e(\vec{a}, \vec{x}, \vec{F})$ by adding the clauses:

1. If $e = \langle 8, d, j \rangle$ and $\lambda x. \phi_d(x, \vec{a}, \vec{f}, \vec{F})$ is total, then $\phi_e(\vec{a}, \vec{f}, \vec{F}) = F_j(\lambda x. \phi_d(x, \vec{a}, \vec{f}, \vec{F}))$.

2. If $e = \langle 9 \rangle$ then $\phi_e(d, \vec{a}, \vec{f}, \vec{F}) = \phi_d(\vec{a}, \vec{f}, \vec{F})$.

This is an example of what is generally called an inductive definition. A large part of generalized computability theory is about inductive definitions and the computational principles that are behind them. In this case, we have to assume that each element in an infinite set of subcomputations will terminate before we accept some computations to terminate. In order to handle this properly, we will have to introduce a generalized concept of computation tree, now a computation tree may be an infinite, well founded tree with countable branching at certain nodes.

It will lead us too far in this introductory course to define all the concepts needed for a mathematically stringent handling of computability relative to functionals. As an appetizer, let us mention two kinds of problems:
1. What are the computable functionals of type 3?

2. What are the functions computable in a fixed functional?

Question 1 is actually hard to answer, there is no obvious characterization of this set. We know that the fan functional, which is defined only for continuous functionals of type 2, is not computable in the sense of Kleene. There are a few other positive and negative results about which functionals that are Kleene-computable, but nothing that is both general and informative. We know more about the answer to the second question, we will address this partly in the next section and partly as a small project discussed in Chapter 5. However, there is no complete and satisfactory characterization of the class of sets of functions that may turn out as the set of functions computable in some functional.

Remark 4.2.3 Scheme ix) may seem a bit dubious, and was actually considered to be a cheat. When we work with Turing machines or with Kleene’s definition of computations relative to functions, the existence of a universal algorithm is an important theorem. However, without scheme ix) we will not be able to prove the existence of a universal algorithm for computations relative to functionals. We will follow Kleene, be pragmatic about it, and claim that including scheme ix) will give us a much more fruitful concept. However, introducing scheme ix) makes the two schemes v) for primitive recursion and vi) for the $\mu$-operator redundant. This, and other facts about computability in functionals, are discussed in the nontrivial Exercise 4.9.

Definition 4.2.4 Let $F$ be a functional of type 2, $f$ a function of type 1. We say that $f$ is Kleene-computable in $F$, $f <_K F$, if for some index $e$ we have that

$$f(x) = \phi_e(x, F)$$

for all $x \in \mathbb{N}$. This definition is extended in the canonical way to the concept $f <_K \vec{f}, \vec{F}$.

We let the 1-section of $F$, $1 - \text{sc}(F)$ be the set

$$1 - \text{sc}(F) = \{ f ; f <_K F \}.$$ 

The next sequence of lemmas should be considered as a small project on which the reader might write an essay:

Lemma 4.2.5 If $f <_K \vec{f}, \vec{F}$ and each $f_i$ in $\vec{f}$ is computable in $\vec{g}, \vec{F}$, then $f <_K \vec{g}, \vec{F}$.

Proof
We may find a primitive recursive function $\rho$ such that if

$$f(x) = \phi_{\rho}(x, \vec{f}, \vec{F})$$

for all $x$ and if

$$f_i(y) = \phi_{\rho_i}(y, \vec{g}, \vec{F})$$
for all $y$ and $i$, then
\[ f(x) = \phi_{\rho(e,e_1,\ldots,e_m)}(\vec{g},\vec{F}) \]
for all $x$.

Lemma 4.2.6 Let $F$ and $G$ be functionals of type 2. Assume that $F(f) = G(f)$ whenever $f \in 1 - sc(F)$. Then $1 - sc(F) = 1 - sc(G)$.

Proof

By induction on the ordinal rank of the computation tree we show that if $\phi_e(\vec{g},G) = a$ then $\phi_e(\vec{g},F) = a$ with the same computation tree.

One of Kleene’s motivations for studying computations relative to functionals was to have a tool for investigating the computational power of quantifiers. Quantification over the natural numbers is captured by the discontinuous functional $^2E$ defined by:

Definition 4.2.7 Let

1. $^2E(f) = 1$ if $\exists a \in \mathbb{N}(f(a) > 0)$.
2. $^2E(f) = 0$ if $\forall a \in \mathbb{N}(f(a) = 0)$.

Definition 4.2.8 Let $F$ and $G$ be functionals of type 2.

a) $F <_K G$ if there is an index $e$ such that $F(f) = \phi_e(G,f)$ for all functions $f : \mathbb{N} \to \mathbb{N}$.

b) $F$ is normal if $^2E <_K F$.

The choice of the term ‘normal’ for these functionals reflects the focus of the early workers of higher type computability. We will investigate computability in $^2E$ more closely in section 4.4. We end this section by showing that the $1$-section of a normal functional will be closed under jumps, see the proof of Lemma 3.3.26 for the definition.

Lemma 4.2.9 Let $F$ be a normal functional. Then $1 - sc(F)$ is closed under jumps.

Proof

Generalizing Kleene’s $T$-predicate, we see that if $g = f'$, there is a computable predicate $T$ such that
\[ g(a) = b \leftrightarrow \exists nT(a,b,n,f). \]

Given $f$ and $a$ we can use $^2E$ to decide if there are $b$ and $n$ such that $T(a,b,n,f)$ holds.

If it does, we may search for the relevant $b$ and output $b + 1$. If it does not, we output 0.
4.3 $2\mathcal{E}$ versus continuity

Let $f_i$ be defined as

\[
\begin{align*}
  f_i(i) &= 1, \\
  f_i(j) &= 0 \text{ if } i \neq j.
\end{align*}
\]

Then $\lim_{i \to \infty} f_i$ is the constant zero function $f$ while $2\mathcal{E}(f) \neq \lim_{n \to \infty} 2\mathcal{E}(f_i)$. This shows that $2\mathcal{E}$ is not continuous. Another way to see this is to observe that $(2\mathcal{E}^{-1}(\{0\}))$ is not an open set, while $\{0\}$ is open in $\mathbb{N}$.

We will now restrict ourselves to considering computations of $\phi_e(\vec{a}, F)$, i.e. computations in one functional argument and some number arguments. We do this in order to save notation, there is no theoretical reason for this restriction.

We will define the $n$’th approximation $\phi^n_e(\vec{a}, F)$ to a computation of $\phi_e(\vec{a}, F)$. There will be four properties to observe

- $\phi^n_e(F, \vec{a})$ will always be defined.
- If $\phi_e(\vec{a}, F) \downarrow$ and $\phi_e(\vec{a}, F) = \lim_{n \to \infty} \phi^n_e(\vec{a}, F)$, then we can tell, in a uniform way, from which $n_0$ the limit is reached.
- If $\phi_e(\vec{a}, F) \downarrow$ and $\phi_e(\vec{a}, F) \neq \lim_{n \to \infty} \phi^n_e(\vec{a}, F)$ we can compute $2\mathcal{E}$ from $F$ in a uniform way.
- If $\phi_e(\vec{a}, F) \downarrow$ we can computably distinguish between the two cases.

We will give most of the details, but the reader should be warned: The rest of this section is an example of how advanced arguments in computability theory might look like. (They often look more advanced than what they really are.)

Definition 4.3.1 We define $\phi^n_e(\vec{a}, F)$ following the cases for the definition of $\phi_e(\vec{a}, F)$ ignoring schemes v), vi) and vii).

i) If $e = \langle 1 \rangle$, we let $\phi^n_e(\vec{a}, F) = \phi_e(\vec{a}, F)$.

ii) If $e = \langle 2, i \rangle$, we let $\phi^n_e(\vec{a}, F) = \phi_e(\vec{a}, F)$.

iii) If $e = \langle 3, q \rangle$, we let $\phi^n_e(\vec{a}, F) = \phi_e(\vec{a}, F)$.

iv) If $e = \langle 4, e', d_1, \ldots, d_m \rangle$, let

\[
\begin{align*}
  \phi^0_e(\vec{a}, F) &= 0, \\
  \phi^0_e(F, \vec{a}) &= \phi_e(\phi^{n_1}_{d_1}(\vec{a}, F), \ldots, \phi^{n_m}_{d_m}(\vec{a}, F)).
\end{align*}
\]

viii) If $e = \langle 8, d, 1 \rangle$, then

\[
\begin{align*}
  \phi^0_e(\vec{a}, F) &= 0, \\
  \phi^0_e(F, \vec{a}) &= F(\lambda x. \phi^0_{d}(x, \vec{a}, F)).
\end{align*}
\]
ix) If \( e = \langle 9 \rangle \), let
\[
\phi^0_e(d, \vec{a}, F) = 0.
\]
\[
\phi^{n+1}_e(d, \vec{a}, F) = \phi^n_e(\vec{a}, F).
\]
\( \phi^n_e(\vec{a}, F) = 0 \) in all other cases.

**Theorem 4.3.2** There are three partial computable functions \( \pi, \eta \) and \( \nu \) defined on the natural numbers such that \( \phi_{\pi(e)}(\vec{a}, F), \phi_{\eta(e)}(\vec{a}, F) \) and \( \phi_{\nu(e)}(\vec{a}, F) \) terminate whenever \( \phi_e(\vec{a}, F) \) terminates, and then

- \( \phi_{\pi(e)}(\vec{a}, F) = 0 \iff \phi_e(\vec{a}, F) = \lim_{n \to \infty} \phi^n_e(\vec{a}, F) \).
- If \( \phi_e(\vec{a}, F) = \lim_{n \to \infty} \phi^n_e(\vec{a}, F) \) then \( \phi^n_e(\vec{a}, F) = \phi_e(\vec{a}, F) \) whenever \( n \geq \phi_{\eta(e)}(\vec{a}, F) \).
- \( \phi_e(\vec{a}, F) \neq \lim_{n \to \infty} \phi^n_e(\vec{a}, F) \) then \( 2E \) is computable in \( F \) via the index \( \phi_{\nu(e)}(\vec{a}, F) \).

**Proof**

We will construct \( \pi, \eta \) and \( \nu \) using the recursion theorem, so we will take the liberty to assume that we know the indices for these functions while defining them.

We will not hide the intuition behind a technically accurate construction of these functions, but to some extent describe in words what to do and why it works. When we explain this, we will assume as an induction hypothesis that our construction works for computations of lower complexity.

We will tell what to do when \( e \) corresponds to one of the schemes we have considered. The ‘otherwise’-case is trivial, since we do not have to prove anything in this case, \( \phi_e \) does not terminate on any input.

If \( e \) corresponds to scheme i), ii) or iii), let \( \pi(e) \) be the index for the constant zero, \( \eta(e) \) the same, and we may without loss of consequences let \( \nu(e) \) also be the same.

Let \( e = \langle 4, e', d_1, \ldots, d_m \rangle \).

Let \( \pi(e) \) be the index for the following enumerated algorithm in \((\vec{a}, F)\):

1. If \( \phi_{\pi(e')}(\vec{a}, F) = \phi_{\pi(d_1)}(\vec{a}, F) = \cdots = \phi_{\pi(d_m)}(\vec{a}, F) = 0 \), let \( \phi_{\pi(e)}(\vec{a}, F) = 0 \) and go to 2., otherwise let \( \phi_{\pi(e)}(\vec{a}, F) = 1 \) and go to 3.

2. Let \( \phi_{\eta(e')}(\vec{a}, F) = 1 + \max\{\phi_{\eta(e')}(\vec{a}, F), \phi_{\eta(d_1)}(\vec{a}, F), \ldots, \phi_{\eta(d_m)}(\vec{a}, F)\} \).

3. Select the least index \( d \in \{e', d_1, \ldots, d_m\} \) such that \( \phi_{\pi(d)}(\vec{a}, F) \neq 0 \). Then use the index \( \phi_{\nu(d)}(F, \vec{a}) \) for \( 2E \) to decide if \( \phi_e(\vec{a}, F) = \lim_{n \to \infty} \phi^n_e(\vec{a}, F) \) or not. If it is, move to 4., otherwise to 5.

4. We let \( \phi_{\pi(e)}(\vec{a}, F) = 0 \). We use \( 2E \) to compute the proper value of \( \phi_{\eta(e)}(\vec{a}, F) \).

5. We let \( \phi_{\pi(e)}(\vec{a}, F) = 1 \), and we let \( \phi_{\nu(e)} = \phi_{\nu(d)} \).
The case ix) is handled in the same way, so we restrict our attention to case viii):
\[
\phi_e(\vec{a}, F) = F(\lambda x. \phi_d(x, \vec{a}, F)).
\]
Let \( f(x) = \phi_e(x, \vec{a}, F) \) and let \( f_n(x) = \phi^n(x, \vec{a}, F) \).

For each \( m \in \mathbb{N} \), define \( g_m(x) \) by the following algorithm:

1. Ask if \( f(x) = \lim_{n \to \infty} f_n(x) \). If yes, continue with 2., otherwise continue with 3.
2. If \( F(f) = F(f_n) \) for all \( n \) such that \( m \leq n \leq \phi_{\eta(d)}(x, F, \vec{a}) \), let \( g_m(x) = f_{\phi_{\eta(d)}(x, \vec{a}, F)}(x) \).
   Otherwise choose the least \( n \geq m \) such that \( F(f) \neq F(f_n) \) and let \( g_m(x) = f_n(x) \).
3. By the induction hypothesis, \( \phi_{\eta(d)}(x, \vec{a}, F) \) provides us with an index to compute \( 2E \) from \( F \). Use \( 2E \) to ask if
   \[
   \exists n \geq m(F(f) \neq F(f_n)).
   \]
   If not, let \( g_m(x) = f(x) \), otherwise let \( g_m(x) = f_n(x) \) for the least such \( n \).

In both cases, we will let \( g_m(x) = f(x) \) if \( F(f) = F(f_n) \) for all \( n \geq m \), while \( g_m(x) = f_n(x) \) for the least counterexample otherwise. In the case 3. this is done explicitly. In case 2. we must use that if \( n > \phi_{\eta(d)}(\vec{a}, F) \), then \( f_n(x) = f_{\phi_{\eta(d)}(x, F, \vec{a})}(x) \), a fact that follows from the induction hypothesis.

Thus for each \( m \) we may computably decide if \( \exists n \geq m(F(f) \neq F(f_n)) \) by
\[
\exists n \geq m(F(f) \neq F(f_n)) \iff F(f) = F(g_m)). \tag{4.2}
\]
From 4.2 we can decide if
\[
\exists n(F(f) \neq F(f_n)).
\]
If not, we let \( \phi_{\pi(e)}(\vec{a}, F) = \phi_{\eta(e)}(\vec{a}, F) = 0 \).

If there exists one such \( n \) we want to decide in a computational way if there are infinitely many of them or not. We may use the same kind of construction as that of \( g_m(x) \) to compute a function \( h(x) \) such that \( h = f \) if \( F(f) \neq F(f_n) \) for infinitely many \( n \), while \( h = f_n \) for the largest \( n \) such that \( F(f) \neq F(f_n) \) otherwise. We must split the construction into the same cases as for the construction of \( g_m(x) \), and we must rely on 4.2 when \( 2E \) is not available. Then we get
\[
\exists n \forall m \geq n(F(f) = F(f_m)) \iff F(f) \neq F(h)). \tag{4.3}
\]
In case both sides of equivalence 4.3 are positive, we can use equivalence 4.2 and the \( \mu \)-operator to find \( \phi_{\eta(e)}(\vec{a}, F) \). It remains to show how to compute \( 2E \) from \( F \) in case both sides of equivalence 4.3 are negative. Actually, we will use the same trick one third time. Assume that \( F(f) \neq F(f_n) \) for infinitely many \( n \). Let \( g \) be given. We want to decide
\[
\exists x(g(x) \neq 0).
\]

We construct \( f'(x) \) as follows:
• If \( f(x) \neq \lim_{n \to \infty} f_n(x) \) use \( 2E \) to decide if \( \exists x (g(x) \neq 0) \). If this is the case, select the least such \( x \), select the least \( n \geq x \) such that \( F(f) \neq F(f_n) \) and let \( f'(x) = f_n(x) \). If it is not the case, let \( f'(x) = f(x) \).

• Otherwise, ask if there is an \( x \leq \phi_{q(d)}(x, \vec{a}, F) \) such that \( g(x) \neq 0 \). In both cases, do as above.

Then
\[
\exists x (g(x) \neq 0) \iff F(f) \neq F(f')
\]

and we are through.

This ends our proof of the theorem.

**Remark 4.3.3** Our use of the recursion theorem is sound, but requires some afterthought by the inexperienced reader. What is required in order to make the definition of \( \pi, \eta \) and \( \nu \) sound is to view them as Turing-computable functions operating on indices for computable functionals of higher types. We have given explicit constructions of \( \phi_{\pi(e)} \) etc. with self reference, and analyzing exactly which combinations of the schemes that are required, we can express \( \pi(e) \) etc. as functions of the Turing indices for \( \pi, \eta \) and \( \nu \). Then, by the recursion theorem, a solution exists.

This theorem has an interesting corollary. A dichotomy theorem is a theorem stating that a general situation splits into two nice, in a sense opposite, cases. One case should then contain more information than just the negation of the other. We have observed that if \( F \) is a normal functional, then \( 1 - \text{sc}(F) \) will be closed under jump. This is actually a characterization, but in a very strong sense. Clearly any 1-section will be closed under the relation ‘computable in’.

We defined a function to be of c.e.degree if it is Turing equivalent with the characteristic function of a c.e. set. This concept may of course be relativized to any function \( g \). Finally, we say that a set \( Y \) of functions is generated from a set \( X \) of functions if \( Y \) consists of all functions computable in some \( g_1, \ldots, g_n \) from \( X \). We then have

**Corollary 4.3.4** Let \( F \) be a functional of type 2. Then one of two will be the case:

1. \( 2E <_K F \), i.e. \( F \) is normal.

2. There is some \( f \in 1 - \text{sc}(F) \) such that \( 1 - \text{sc}(F) \) is generated by the elements in \( 1 - \text{sc}(F) \) of c.e.(f)-degree.

**Proof**

If there is one terminating computation \( \phi_e(\vec{a}, F) \) such that
\[
\lim_{n \to \infty} \phi_e^n(\vec{a}, F) \neq \phi_e(\vec{a}, F)
\]

then \( F \) is normal by Theorem 4.3.2.

Otherwise, let \( f(\langle e, n, \vec{a} \rangle) = \phi_e^n(\vec{a}, F) \). Clearly, \( f \in 1 - \text{sc}(F) \). Let \( g \in 1 - \text{sc}(F) \)
be given, \( g(x) = \phi_v(x,F) \) for all \( x \). Then \( g \) is computable in \( f \) and the c.e.(f) set
\[
A = \{(x,n) : \exists m \geq n (\phi^m_v(x,F) \neq \phi^n_v(x,F))\}
\]
which again is computable in \( F \) via \( \lambda x.\phi_{\eta(x)}(x) \).

Kleene computability is extended to functionals of all finite types, and this dichotomy actually still holds. However, working with functionals in which \( 2^E \) is not computable, one may as well work with what is known as the hereditarily continuous functionals. We will touch a little bit on this in section 4.5.

### 4.4 The Hyperarithmetical sets

**Definition 4.4.1** A set is hyperarithmetical if it is computable in \( 2^E \).

We have cheated a bit, and used a characterization due to Kleene as our definition. For a more systematic introduction to hyperarithmetical sets and higher computability theory in general, see Sacks [3].

The term ‘hyperarithmetical’ indicates that this is an extension of the arithmetical sets in a natural way, and this is exactly what the intention is. The arithmetical sets will be the sets that can be defined by first order formulas in number theory. For each Gödel number \( \alpha \) of a first order formula \( \psi(x) \) with one free variable, we can, in a primitive recursive way, find an index \( \pi(\alpha) \) such that
\[
\lambda x.\phi_{\pi(\alpha)}(x,2^E)
\]
will be the characteristic function of
\[
\{x ; \psi(x)\}.
\]
Using the enumeration scheme ix) we can then find a subset \( B \) of \( \mathbb{N} \times \mathbb{N} \) that is computable in \( 2^E \) and such that each arithmetical set \( A \) is a section \( B_x = \{y ; (x,y) \in B\} \) of \( B \). Then every set arithmetical in \( B \) will also be hyperarithmetical, and so forth. In a sense, we may say that \( 2^E \) provides the ‘arithmetical’ while the enumeration scheme provides the ‘hyper’.

### 4.4.1 Trees

If we want to analyze computations relative to functionals more closely, we need a precise notion of a computation tree. Although there is a common intuition behind all our concepts of trees, various needs require various conventions. For the rest of this section we will use

**Definition 4.4.2**

a) A tree will be a nonempty set \( T \) of sequences of natural numbers, identified with their sequence numbers, that is closed under initial segments. The elements of \( T \) will be called nodes.

b) \( \langle \rangle \), i.e. the empty sequence, is the root node of the tree.
c) A leaf node of a tree $T$ will be a node in $T$ with no proper extension in $T$.

d) A branch in a tree $T$ will be a maximal, totally ordered subset of $T$. If a branch is infinite, we identify it with the corresponding function $f : \mathbb{N} \rightarrow \mathbb{N}$ and if a branch is finite, we identify it with its maximal element.

e) A decorated tree will be a tree $T$ together with a decoration, i.e. a map $f : T \rightarrow \mathbb{N}$.

f) If $T_i$ is a tree for each $i \in I$, where $I \subseteq \mathbb{N}$, we let $\langle T_i \rangle _{i \in I}$ be the tree $T$ consisting of the empty sequence together with all sequences $i * \sigma$ such that $i \in I$ and $\sigma \in T_i$.

g) If $a \in \mathbb{N}$ and $(T_i, f_i)$ are decorated trees for each $i \in I$, we let $\langle a, (T_i, f_i) \rangle _{i \in I}$ be the decorated tree $(T, f)$ defined by

- $T = \langle T_i \rangle _{i \in I}$
- $f(\langle \rangle ) = a$
- $f(i * \sigma) = f_i(\sigma)$ when $i \in I$ and $\sigma \in T_i$.

A computation tree will be a decorated tree. Since it actually is the decorations that will be of interest, we sometimes take the liberty to identify them with the corresponding nodes with decoration in the tree. This liberty will be visible when we discuss the root and the leaves of a computation tree.

**Definition 4.4.3** Let $\phi_e^e(2E, \vec{a}) = b$. We define the computation tree of the computation by recursion as follows

i) $e = \langle 1 \rangle$: Let the root node also be the only leaf node and be decorated with $(e, \vec{a}, b)$.

ii) $e = \langle 2, i \rangle$. Act as in case i).

iii) $e = \langle 3, q \rangle$. Act as in case i).

iv) $e = \langle e', d_1, \ldots, d_m \rangle$: Let $(T_1, f_1), \ldots, (T_m, f_m)$ be the computation trees of $\phi_{d_i}(2E, \vec{a}) = c_i$ resp. and let $(T_{m+1}, f_{m+1})$ be the computation tree of $\phi_{e'}(2E, c_1, \ldots, c_m) = b$.

Let the computation tree of $\phi_e(2E, \vec{a}) = b$ be

$$\langle \langle e, \vec{a}, b \rangle, T_i \rangle _{i=1}^{m+1}.$$ 

vii) $e = \langle 8, d \rangle$: For each $i \in \mathbb{N}$, let $(T_i, f_i)$ be the computation tree of $\phi_d(i, \vec{a}, 2E)$ and let the computation tree of $\phi_e(\vec{a}, 2E) = b$ be

$$\langle \langle e, \vec{a}, b \rangle, (T_i, f_i) \rangle _{i \in \mathbb{N}}.$$ 

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The rest of the cases are similar, and the details are left for the reader.

**Lemma 4.4.4** Whenever $\phi_e(\bar{a},2E)$ terminates, then the computation tree will be computable in $2E$.

**Proof**
We see how the computation tree is constructed from the computation trees of the immediate subcomputations. This construction is clearly computable in $2E$, and using the recursion theorem for computations relative to $2E$ we get a uniform algorithm for computing the computation tree from a computation.

Recall that a tree is well founded if there is no infinite branch in the tree. It is easy to see that the computation trees of terminating computations will be well founded.

**Definition 4.4.5** A pre-computation tree will be a decorated tree that locally looks like a computation tree, i.e. each node will be of the form $\langle d, \bar{a}, b \rangle$, the leaf nodes will correspond to initial computations, and other nodes relate to their immediate subnodes as in Definition 4.4.3.

**Lemma 4.4.6**

a) The concept of a pre-computation tree is arithmetical.

b) A pre-computation tree that is well founded is actually a computation tree.

c) If $\phi_e(\bar{a},2E)$ terminates, there is exactly one pre-computation tree with a root node on the form $\langle e, \bar{a}, b \rangle$, and then $b$ will be the value of $\phi_e(\bar{a},2E)$.

**Proof**
a) is trivial. In order to prove b) we observe that a subtree of a well founded pre-computation tree is itself a well founded pre-computation tree. Since facts may be proved by induction on the subtree ordering of well founded trees, we may use this kind of induction to prove the statement. The details are trivial.

The set of computations relative to $2E$ was defined by induction, and then we may prove facts by induction over this construction. c) is proved this way in a trivial manner.

### 4.4.2 $\Pi^0_k$-sets etc.

In this section we will let ‘computable’ mean ‘Turing-computable’.

**Definition 4.4.7**

a) A product set will be any product of the sets $\mathbb{N}$ and $\mathbb{N}^\mathbb{N}$ in any finite number and order. We define the arithmetical hierarchy of subsets of product sets $X$ as follows:

i) A set $A$ is a $\Sigma^0_0$-set and a $\Pi^0_0$-set if $A$ is computable. (In the literature you may find that $A$ is supposed to be definable using only bounded quantifiers in order to be in these classes. This distinction does not matter for our applications.)
ii) A set $A \subset X$ is $\Sigma^0_{k+1}$ if for some $\Pi^0_k$ subset $B$ of $\mathbb{N} \times X$ we have that
\[
\bar{x} \in A \iff \exists x \in \mathbb{N}((x, \bar{x}) \in B).
\]

iii) A set $A \subset X$ is $\Pi^0_{k+1}$ if for some $\Sigma^0_k$ subset $B$ of $\mathbb{N} \times X$ we have that
\[
\bar{x} \in A \iff \forall x \in \mathbb{N}((x, \bar{x}) \in B).
\]

iv) $A$ is a $\Delta^0_k$-set if $A$ is both $\Pi^0_k$ and $\Sigma^0_k$.

b) We define the \textit{analytical hierarchy} in the same fashion, but will only need the first level here:

i) A set $A \subseteq X$ is $\Pi^1_1$ if there is some arithmetical set $B \subseteq \mathbb{N}^\mathbb{N} \times X$ such that
\[
\bar{x} \in A \iff \forall f \in \mathbb{N}^\mathbb{N}((f, \bar{x}) \in B).
\]

ii) A set $A \subseteq X$ is $\Sigma^1_1$ if there is some arithmetical set $B \subseteq \mathbb{N}^\mathbb{N} \times X$ such that
\[
\bar{x} \in A \iff \exists f \in \mathbb{N}^\mathbb{N}((f, \bar{x}) \in B).
\]

iii) $A$ is $\Delta^1_1$ if $A$ is both $\Sigma^1_1$ and $\Pi^1_1$.

\textbf{Lemma 4.4.8} Let $X$ be a product set, $A \subseteq X$. Then $A$ is $\Sigma^0_1$ if and only if $A$ is the domain of a partial computable function of arity $X$.

\textit{Proof}\:
For product sets where all factors are $\mathbb{N}$, these are two of the characterizations of c.e.-sets. The proof of the equivalence generalizes trivially to all product sets.

In section 4.1 we showed how each product class either is homeomorphic to $\mathbb{N}$ or to $\mathbb{N}^\mathbb{N}$. We used pairing-functions $(\cdot, \cdot)$ for pairs of numbers, pairs of functions or for a pair consisting of one number and one function for this. These functions can be used to show that multiple quantifiers of the same kind can be reduced to one, and that number quantifiers can be ‘eaten up’ by function quantifiers. For a complete proof, we need one more homeomorphism:

\textbf{Lemma 4.4.9} $(\mathbb{N}^\mathbb{N})^\mathbb{N}$ is homeomorphic to $\mathbb{N}^\mathbb{N}$.

\textit{Proof}\:
We use the observation $(\mathbb{N}^\mathbb{N})^\mathbb{N} \cong \mathbb{N}^{\mathbb{N} \times \mathbb{N}} \cong \mathbb{N}^\mathbb{N}$ since $\mathbb{N} \times \mathbb{N} \cong \mathbb{N}$. Precisely, if $\{f_i\}_{i \in \mathbb{N}}$ is a sequence of functions, we code it as one function by
\[
(f_i)_{i \in \mathbb{N}}(\langle n, m \rangle) = f_n(m).
\]
This is a computable bijection.

All these coding functions will have inverses for each coordinate. We use $(\cdot)_i$, for the inverse at coordinate $i$, letting pairing functions have coordinates 1 and 2. We can then perform the following reductions of quantifiers. Note that all the dual reductions will also hold.
Lemma 4.4.10  For each relation $R$ the following equivalences hold:

i) $\exists n \exists m R(n, m) \Leftrightarrow \exists k R((k)_1, (k)_2)$.

ii) $\exists f \exists g R(f, g) \Leftrightarrow \exists h R((h)_1, (h)_2)$.

iii) $\exists f \exists n R(f, n) \Leftrightarrow \exists g R((g)_1, (g)_2)$ (Where the decoding is different from the one in ii).

iv) $\forall n \exists m R(n, m) \Leftrightarrow \exists f \forall n R(n, f(n))$.

v) $\forall n \exists f R(n, f) \Leftrightarrow \exists g \forall n R(n, g(n))$.

The proofs are trivial and are left for the reader as Exercise 4.10.

Lemma 4.4.11  

a) The classes $\Sigma^0_k$, $\Pi^0_k$, $\Sigma^1_1$ and $\Pi^1_1$ are closed under finite unions and finite intersections. Moreover, the $\Sigma$-classes are closed under $\exists n \in \mathbb{N}$ and the $\Pi$-classes are closed under $\forall n \in \mathbb{N}$.

b) The classes $\Pi^1_1$ and $\Sigma^1_1$ are closed under number quantifiers.

c) $\Pi^1_1$-normal form theorem

If $A \subseteq X$ is $\Pi^1_1$, there is a computable set $R \subseteq \mathbb{N}^2 \times X$ such that

$$\vec{x} \in A \Leftrightarrow \forall f \exists n R(f, n, \vec{x}).$$

In the proof we use standard prenex operations and the reductions described in Lemma 4.4.10 The details are left for the reader, see Exercise 4.10.

Theorem 4.4.12  Let $\Gamma$ be one of the classes $\Pi^0_k$, $\Sigma^0_k$ ($k \geq 1$), $\Pi^1_k$ or $\Sigma^1_k$ (or $\Pi^1_1$, $\Sigma^1_1$ as defined in Exercise 4.11). For each product set $X$ there is a universal $\Gamma$-set $A$ in $\mathbb{N} \times X$, i.e. such that for each $\Gamma$-set $B \subseteq X$ there is an $n \in \mathbb{N}$ such that for all $\vec{x} \in X$:

$$\vec{x} \in B \Leftrightarrow (n, \vec{x}) \in A.$$ 

Proof

Let $\tilde{\Gamma}$ be the dual of $\Gamma$, i.e. the set of complements of sets in $\Gamma$. Clearly, if the property holds for $\Gamma$, it also holds for $\tilde{\Gamma}$. Since the $\Pi$- and $\Sigma$-classes are duals of each other, it is sufficient to prove the lemma for one of each pair.

The existence of universal algorithms together with Lemma 4.4.8 ensures the theorem to hold for $\Sigma^0_k$. Then any class $\Gamma$ of sets definable from $\Sigma^0_k$-sets using a fixed quantifier prefix will have universal sets for each product set. Each of our classes $\Gamma$ is either one of these or the dual to one of these. For $\Pi^1_1$ we need the normal form theorem, see Lemma 4.4.11 c).

We will end this subsection by viewing the connection between $\Pi^1_1$-sets and well founded trees.

Let $A \subseteq \mathbb{N}$ be a $\Pi^1_1$-set written in its normal form

$$m \in A \Leftrightarrow \forall f \exists n R(m, n, f).$$
Since \( R \) is computable, there is, by the finite use property, a computable relation \( R^+ \) on \( \mathbb{N}^3 \) such that
\[
m \in A \iff \forall f \exists n \exists k R^+(m,n,f(k)).
\]
For each \( m \) we let \( T_m \) be the tree of finite sequences \( \sigma \) such that
\[
\forall n \leq \text{lth}(\sigma) \forall \tau \prec \sigma \neg R^+(m,n,\tau).
\]
Then \( T_m \) will be a tree, and we will have
\[
m \in A \iff T_m \text{ is well founded}.
\]
If we put a bit more effort into the arguments, we can actually show that we may choose \( R^+ \) to be primitive recursive, this is connected with the relativized Kleene’s \( T \)-predicate, and thus each \( \Pi_1 \) subset of \( \mathbb{N} \) is \( m \)-reducible to the set of indices for well founded primitive recursive trees. On the other hand, this set is easily seen to be \( \Pi_1 \) itself. Thus the class of \( \Pi_1 \)-sets contains an \( m \)-maximal element, just like the c.e. sets.

In the next subsection we will see more analogies between the c.e. sets and the \( \Pi_1 \)-sets.

### 4.4.3 Semicomputability in \( ^2E \) and Gandy Selection

**Definition 4.4.13** A subset \( A \subseteq \mathbb{N} \) is **semicomputable in \( ^2E \)** if there is an index \( e \) such that for all \( a \in \mathbb{N} \):
\[
\phi_e(a, ^2E) \downarrow \iff a \in A,
\]
where \( \downarrow \) still means ‘terminates’.

**Lemma 4.4.14** All \( \Pi_1 \) subsets of \( \mathbb{N} \) are semicomputable in \( ^2E \).

**Proof**

Clearly the set of indices for computable trees is computable in \( ^2E \), so it is sufficient to show that the concept of a well founded tree is semicomputable.

This can be done with a careful use of the recursion theorem: Let \( e \) be any index. Using the enumeration scheme and \( ^2E \) we can define a computable function \( \Phi(e,T) \) such that

- \( \Phi(e,T, ^2E) = 0 \) if \( T \) consists of exactly the empty sequence.
- \( \Phi(e,T, ^2E) = ^2E(\lambda n \phi_e(T_n)) \) if \( T \) is a more complex tree, where \( T_n \) is the set of sequences \( \sigma \) such that \( n * \sigma \in T \).

By the recursion theorem for \( ^2E \), there is an index \( e_0 \) such that \( \Phi(e_0,T, ^2E) = \phi_{e_0}(T, ^2E) \) for all \( T \).

By induction on the well founded trees it is easy to see that \( \phi_{e_0}(T) \) will terminate whenever \( T \) is well founded. In order to prove the other direction, we must inspect the proof of the recursion theorem and see that in this particular case, if \( \phi_{e_0}(T) \) terminates, then either \( T \) consists of the empty sequence only or \( \phi_{e_0}(T_n) \)
must terminate for each \( n \) and be a subcomputation. This is because we only used \( e \) as an index in connection with the enumeration scheme in defining \( \Phi \). It follows that \( T \) must be well founded.

**Theorem 4.4.15** Let \( A \subseteq \mathbb{N} \). Then the following are equivalent:

i) \( A \) is \( \Pi^1_1 \).

ii) \( A \) is semicomputable in \( 2^E \).

**Proof**

Lemma 4.4.14 gives us one direction.

The set \( C \) of computation tuples \( \langle e, \bar{a}, b \rangle \) such that \( \phi_e(2^E, \bar{a}) = b \) is defined by a *positive induction*, i.e. there is a formula \( \Phi(x, X) \) such that atomic subformulas \( t(x) \in X \) will occur positively and such that \( C \) is the least set such that

\[
C = \{ a ; \Phi(a, C) \}.
\]

Then

\[
e \in C \iff \forall B (B \subseteq \{ a ; \Phi(a, B) \} \rightarrow e \in B).
\]

**Remark 4.4.16** This proof is trivial provided the reader has been through a quick introduction to inductive definitions.

In Exercise 5.9 we offer a self-service introduction.

Now is the time to assume that the reader is familiar with well orderings and preferably with ordinal numbers. For readers unfamiliar with this, we offer Exercise 5.11, a guided self-service introduction to the topic.

**Definition 4.4.17** Let \( T \) be a well founded tree on \( \mathbb{N} \).

a) The rank \( |T| \) of \( T \), is the ordinal rank of \( T \) seen as a well founded relation, where \( \sigma \leq \tau \iff \tau \prec \sigma \).

b) The rank \( |\sigma|_T \) of \( \sigma \in T \), is the value of the rank function for \( T \) on \( \sigma \).

c) If \( T \) is not well founded, we let \( |T| = \infty \) considered to be larger than any ordinal number.

**Theorem 4.4.18** There is a two-place function \( \Phi \) partially computable in \( 2^E \) such that \( \phi(T, S) \) terminates exactly when both \( S \) and \( T \) are trees, and at least one of them is well founded, and then

1. If \( |T| \leq |S| \) then \( \Phi(T, S) = 0 \)

2. If \( |S| < |T| \) then \( \Phi(T, S) = 1 \).

**Proof**

We will use the recursion theorem for \( 2^E \), and the argument is an elaboration on the argument showing that the set of well founded trees is semicomputable in \( 2^E \).
We define $\Phi(S,T)$ by cases, using self reference, and the solution using the construction behind the recursion theorem will give us the result. The set of trees is computable in $2E$, so for the sake of convenience, we assume that both $T$ and $S$ are trees. Then

- If $T$ consists of only the empty sequence, let $\Phi(T,S) = 0$.
- If $S$ consists of only the empty sequence, but $T$ contains more, let $\Phi(T,S) = 1$.
- If neither $T$ nor $S$ contains just the empty sequence, let

$$\Phi(S,T) = 0 \iff \forall n \exists m (\Phi(T_n,S_m) = 0),$$

where we use $2E$ to decide this, and let 1 be the alternative value.

It is easy to see by induction on the rank that if one of the trees is well founded, then $\Phi$ does what it is supposed to do. If neither $T$ nor $S$ are well founded, $\Phi(T,S)$ will not terminate. This is however not important, and we skip the argument.

Since every terminating computation is associated with a well founded computation tree, we may think of the rank of the computation tree as a measure of the length of the computation.

**Definition 4.4.19** Let $\phi_e(2E,\vec{a}) = b$. The length $|\langle e, \vec{a}, b \rangle|$ of the computation will be the ordinal rank of the corresponding computation tree. If $\langle e, \vec{a}, b \rangle$ is not the tuple of a terminating computation, we let $|\langle e, \vec{a}, b \rangle| = \infty$.

Theorem 4.4.18 has an interesting application, the *Stage Comparison Theorem*:

**Corollary 4.4.20** There is a function $\Psi$ partially computable in $2E$ such that whenever $\sigma = \langle e, \vec{a}, b \rangle$ and $\sigma' = \langle e', \vec{a'}, b' \rangle$ are tuples, then $\Psi(\sigma, \sigma')$ terminates if and only if at least one of $\sigma$ and $\sigma'$ is a genuine computation tuple, and then

$$\Psi(\sigma, \sigma') = 0 \iff |\sigma| \leq |\sigma'|.$$

We have proved a selection theorem for ordinary c.e. sets, the selection was performed by searching for the least pair where the second element was a witness to the fact that the first element was in the c.e. set in question. Gandy showed how we may combine this idea and the ordinal length of computations to prove a selection theorem for computations in $2E$:

**Theorem 4.4.21 (Gandy Selection)**

Let $A \subseteq \mathbb{N} \times \mathbb{N}$ be c.e. in $2E$. Then there is a selection function for $A$ computable in $2E$.

**Proof**

We will use the recursion theorem for $2E$. Let

$$A = \{ (a,b) : \phi_{e_0}(a,b,2E) \}. $$

Let $f$ be computable in $2E$ satisfying
1. \( f(e, a, k) = 0 \) if \(|\langle e_0, a, k \rangle| < |\langle e, a, k + 1 \rangle|\).

2. \( f(e, a, k) = f(e, a, k + 1) + 1 \) if \( \phi(e, a, k + 1, 2E) \downarrow \) and \(|\langle e, a, k + 1 \rangle| \leq |\langle e_0, a, k \rangle|\).

By the recursion theorem there is an index \( e_1 \) such that

\[ f(e_1, a, k) = \phi(e_1, a, k, 2E). \]

We claim that \( \lambda a.\phi(e_1, a, 0, 2E) \) is a selection function for \( A \). Let \( a \in \mathbb{N} \) be given.

First observe that if \( \phi(e_0, a, k, 2E) \downarrow \) then \( f(e, a, k) \) will terminate, so \( \phi(e_0, a, k, 2E) \downarrow \).

Moreover, observe that if \( \phi(e_1, a, k', 2E) \downarrow \) whenever \( k' < k \).

Thus

\[ \exists k \phi(e_0, a, k, 2E) \downarrow \Rightarrow \phi(e_1, a, k, 2E) \downarrow. \]

Now assume that \( \phi(e_1, a, 0, 2E) \downarrow \). If we look at the computation of \( \phi(e_1, a, 0, 2E) \) we see that we will have \( \phi(e_1, a, 1, 2E), \phi(e_1, a, 2, 2E), \ldots \) as subcomputations as long as part 2 of the algorithm for \( f \) is followed. The ranks of these computations will be a decreasing sequence of ordinals, and this sequence must come to an end. It comes to an end exactly when we hit a \( k_0 \) such that

\[ |\langle e_0, a, k_0 \rangle| < |\langle e_1, a, k_0 + 1 \rangle|. \]

Then \( (a, k_0) \in A \) and backtracking the value of \( \phi(e_1, a, k', 2E) \) for \( k' \leq k_0 \) we see that \( \phi(e_1, a, k', 2E) = k_0 - k' \).

Consequently \( \lambda a.\phi(e_1, a, 0, 2E) \) will be a selection function for \( A \)

**Remark 4.4.22** This proof is of course uniform in \( e_0 \).

**Corollary 4.4.23** Let \( A \subseteq \mathbb{N} \). Then the following are equivalent:

i) \( A \) is computable in \( 2E \).

ii) Both \( A \) and \( \mathbb{N} \setminus A \) are semicomputable in \( 2E \).

**Proof**

Clearly, if \( A \) is computable in \( 2E \), then both \( A \) and its complement will be semicomputable in \( 2E \).

Now assume that both \( A \) and \( \mathbb{N} \setminus A \) are semicomputable in \( 2E \).

Let

\[ B = \{(0, n) ; n \in A\} \cup \{(1, m) ; m \notin A\}. \]

\( B \) is semicomputable, and \( B \) is the graph of a function. By Gandy selection this function must be computable in \( 2E \). The result follows.

**4.4.4 Characterising the hyperarithmetical sets**

There is almost nothing left for us to do in this subsection. We have shown that a set \( A \) is semicomputable in \( 2E \) if and only if it is \( \Pi^1_1 \). We have shown that a set \( B \) is computable in \( 2E \) if both \( B \) and its complement are semicomputable in \( 2E \). This gives us
Corollary 4.4.24 A set is hyperarithmetical if and only if it is $\Delta^1_1$.

Remark 4.4.25 Our proof of Corollary 4.4.24 is flavored by computability theory, but there are alternative proofs in the literature. The result, in a different form, goes back to Suslin in 1917. He essentially showed that a set $A \subseteq \mathbb{N}^\mathbb{N}$ is Borel if and only if both $A$ and its complement are projections of closed sets in $(\mathbb{N}^\mathbb{N})^2$, i.e. that relativized $\Delta^1_1$ is the same as Borel.

There are numerous analogies between the pair (Computable, computably enumerable) and the pair $(\Delta^1_1, \Pi^1_1)$. Some of these are left for the reader as Exercise 4.12.

There is a close connection between hyperarithmetical theory and fragments of set theory. For readers familiar with axiomatic set theory, we offer a self-service introduction in Exercise 5.12.

4.5 Typed \(\lambda\)-calculus and PCF

Kleene’s definition of computations relative to functionals can be extended to functionals of even higher types than two. This will, however, be a too specialized topic to be introduced in this text. Kleene’s concept turned out to be useful in definability theory, the theory of what may be defined with the help of certain principles, but we have moved quite a bit away from what we might call “genuine computability”.

The concept of higher type computability is nevertheless of interest also when genuine computability is the issue, e.g. in theoretical computer science. In this section we will give a brief introduction to PCF. PCF is a formal programming language for computing with typed objects. It has its roots in work by Platek, Scott developed the first version of it as a formal logic for computability and Plotkin gave it the form we will be investigating. The intuition should be that we are operating with hereditarily partial, monotone and continuous functionals. We will explain this better when needed.

4.5.1 Syntax of PCF

Definition 4.5.1 We define the formal types, or just types as a set of terms for types as follows:

1. $\iota$ and $o$ are types. These are called the formal base types.
2. If $\sigma$ and $\tau$ are types, then $(\sigma \rightarrow \tau)$ is a type.

Remark 4.5.2 When we give a semantical interpretation of PCF, we will associate a mathematical object to each formal type. We think of $\iota$ as denoting the set of natural numbers, $o$ as denoting the set of boolean values and $(\sigma \rightarrow \tau)$ as denoting the relevant set of functions mapping objects of type $\sigma$ to objects of type $\tau$. Since we will use these objects to interpret algorithms, we have to
interpret nonterminating algorithms as well. We will do so by including an element for 'the undefined' in the base types, and carry this with us for higher types.

As usual the term language will consist of variables, constants and combined terms. What is new is that each term will be typed, that some constants are treated as function symbols, and that we need type matching when forming combined terms.

Definition 4.5.3 The terms in PCF are inductively defined as follows:

1. Variables:
   For each type \( \sigma \) there is an infinite list \( x_1^\sigma \) of variables of type \( \sigma \).

2. Constants:
   In PCF we have the following typed constants:
   - \( k_n \) of type \( \iota \) for each natural number \( n \).
   - \( tt \) and \( ff \) of type \( o \).
   - \( Z \) of type \( \iota \rightarrow o \).
   - \( S \) and \( P \) of type \( \iota \rightarrow \iota \).
   - \( \supset \) and \( \supsetendo \) of types \( o \rightarrow (\iota \rightarrow (\iota \rightarrow \iota)) \) and \( o \rightarrow (o \rightarrow (o \rightarrow o)) \) resp.
   - \( Y_\sigma \) of type \( (\sigma \rightarrow \sigma) \rightarrow \sigma \) for each type \( \sigma \).

3. Combined terms:
   - If \( M \) is a term of type \( \sigma \rightarrow \tau \) and \( N \) is a term of type \( \sigma \), then \( (MN) \) is a term of type \( \tau \). This construction is called application.
   - If \( M \) is a term of type \( \tau \) and \( x \) is a variable of type \( \sigma \), then \( (\lambda x.M) \) is a term of type \( \sigma \rightarrow \tau \). This construction is called abstraction.

The intuition behind these constants and terms are as follows:
0 will denote the zero-element in \( \mathbb{N} \), \( tt \) and \( ff \) the two truth values.
\( Z \) will test if a number is zero or not.
\( S \) and \( P \) are the successor and predecessor operators on the natural numbers.
\( \supset \) will select the second or third argument depending on the boolean value of the first argument. Thus there will be one for each base type.
\( Y_\sigma \) will denote the fixed point operator. The idea is that whenever \( f : \sigma \rightarrow \sigma \), then \( f \) will have a least fixed point \( a \) of type \( \sigma \). Our challenge will be to formalize this in the formal theory and to find mathematical interpretations of these types such that this makes sense.
Application \( (MN) \) simply mean that \( M \) is thought of as a function, \( N \) as an argument and \( (MN) \) then just denotes the application. It may be confusing that we are not using the more conventional notation \( M(N) \), but this has turned out to be less convenient in this context.
If $M$ is a term of type $\tau$, $M$ may contain a variable $x$ of a type $\sigma$. The intuition is that $M$ denotes an unspecified object of type $\tau$, an object that becomes specified when we specify the value of $x$. $(\lambda x. M)$ will denote the function that maps the specification of $x$ to the corresponding specification of $M$.

The reader may have noticed that we have omitted a few parentheses here and there. We will do so systematically, both for types and terms.

If $f$ is a function of type $(\sigma \to (\tau \to \delta))$ we will view $f$ as a function of two variables of types $\sigma$ and $\tau$. We will write $\sigma, \tau \to \delta$ for such types. If then $M$ is of type $\sigma, \tau \to \delta$, $N$ is of type $\sigma$ and $K$ is of type $\tau$ we will write $MNK$ instead of $(((M)N)(K))$. This means that the application is taken from left to right.

**Lemma 4.5.4** Each type will be of the form $\sigma = \tau_1, \ldots, \tau_n \to b$ where $n \geq 0$ and $b$ is one of the base types.

The proof is easy by induction on the length of $\sigma$ seen as a word in an alphabet.

With this convention we see that the type of $\supset_i$ is $o, i, i \to i$ and the type of $\supset_o$ is $o, o, o \to o$. We then view these as functions of three variables.

### 4.5.2 Operational semantics for PCF

An operational semantics for a language designed to describe algorithms will be the specification of how to carry out step-by-step calculations or computations. The operational semantics for PCF will be a set of rules for how to rewrite terms in order to ‘compute’ the value. This is of course most relevant when we are dealing with combined terms of base types.

In PCF we have a similar distinction between free and bounded occurrences of variables as in first order logic, $x$ becomes bounded in $\lambda x. M$. As for first order languages, one term $N$ may be substituted for a variable $x$ in a term $M$ if no variables free in $N$ becomes bounded after the substitution. We write $M_K$ for the result of the substitution, always assuming that $N$ is substitutable for $x$ in $M$.

**Definition 4.5.5** We define the relation $\longrightarrow$, denoted by a long arrow, as the reflexive and transitive closure of the following one-step reductions (we assume that the typing is correct):

- $Zk_0 \longrightarrow tt$.
- $Zk_{n+1} \longrightarrow ff$.
- $Pk_{n+1} \longrightarrow k_n$.
- $Sk_n \longrightarrow k_{n+1}$.
- $\supset_b ttMN \longrightarrow M$ where $b$ is a base type.
$\Gamma b. f f M N \rightarrow N$ where $b$ is a base type.

$(\lambda x. M) N \rightarrow M \tilde{x}$.

$Y_\sigma M \rightarrow M (Y_\sigma M)$.

$M \rightarrow M' \Rightarrow MN \rightarrow M' N$.

$N \rightarrow N' \Rightarrow MN \rightarrow MN'$.

$\lambda x. M \rightarrow \lambda y. M x y$.

The last item tells us that we may replace a bounded quantifier with another one not occurring in $M$. As an example we will see how we may compute $f(x) = 2x$ on the natural numbers:

**Example 4.5.6** The function $f(x) = 2x$ is defined by primitive recursion from the successor operator as follows:

- $f(0) = 0$
- $f(S(x)) = S(S(f(x)))$.

For the sake of readability, let $g$ be a variable of type $\iota \rightarrow \iota$. Consider the term $M$ of type $\iota \rightarrow \iota$ defined by

$$M = Y_{\iota \rightarrow \iota} \lambda g. \lambda x^\iota (\supset_{\iota} (Zx) k_0 S(S(g(Px))))$$

If we look at the expression $\lambda g. \lambda x^\iota (\supset_{\iota} (Zx) k_0 S(S(g(Px))))$ we see that to each $g$ we define a function that takes the value 0 on input 0 and $g(x - 1) + 2$ for positive inputs $x$. The least fixed point of this operator is exactly the function $f$. The verification of e.g. $M k_2 \rightarrow k_4$ is a lengthy process.

**Example 4.5.7** We will show how the $\mu$-operator can be handled by PCF, i.e. we will define a term of type $(\iota \rightarrow \iota) \rightarrow \iota$ that must be interpreted as the $\mu$-operator.

Note that $F(f) = \mu x. f(x) = 0$ can in the same sense be defined in the following way: $F(f) = G(f, 0)$ where

$$G(f, k) = 0 \text{ if } f(k) = 0.$$  

$$G(f, k) = G(f, k + 1) + 1 \text{ if } f(k) > 0.$$  

$G$ will be the fixed point of a PCF-definable operator, and then the $\mu$-operator is definable.

The conditionals $\supset_{\iota}$ and $\supset_{\sigma}$ can be extended to conditionals $\supset_{\sigma}$ for all types $\sigma$, see Exercise 4.13

**Definition 4.5.8** Let $f : \mathbb{N} \rightarrow \mathbb{N}$ be a partial function. We say that $f$ is PCF-definable if there is a term $M$ of type $\iota \rightarrow \iota$ such that the following are equivalent for all numbers $n$ and $m$:
1. $Mk_n \rightarrow k_m$

2. $f(k)$ is defined and $f(n) = m$.

We used the expression ‘$PCF$-definable’ also for the $\mu$-operator, trusting that the reader accepts this yet not clarified terminology. We then have

**Theorem 4.5.9** The $PCF$-definable functions are exactly the partial computable ones.

**Proof**
We have shown how to handle the $\mu$-operator and indicated how to handle primitive recursion. The full details are then easy and are left for the reader as Exercise 4.14.

**4.5.3 A denotational semantics for $PCF$**

The operational semantics for $PCF$ explains how we may rewrite a term in such a way that we eventually may read off a number or a boolean value as the interpretation of the term. When we give a denotational semantics, we will interpret each term as a mathematical object in some set. One important aspect will be that the reductions of the operational semantics should not alter the denotational interpretations. If we look at the term $M$ for the $f(x) = 2x$ function, we obtain that $Mk_2 \rightarrow k_4$ in the operational semantics, while we want $[Mk_2] = 4$, where we will use $[\cdot]$ for the denotational semantics.

For the rest of this section, we leave out essentially all details. It is not expected that the reader will be able to fill in the details without consulting a textbook on domain theory or some other introduction to a similar topic. Most lemmas etc. will be left without proof.

Consider the term $N = Yιλx.S(x)$. If we try to evaluate $Nk_0$ we see that we get an infinite reduction sequence. Thus the only sensible way to interpret this term is by accepting ‘undefined’ as a possible value, and then using it in this case.

**Definition 4.5.10** Let $\mathbb{N}_\perp$ and $\mathbb{B}_\perp$ be the set $\mathbb{N}$ of natural numbers and the set $\mathbb{B}$ of boolean values $\{true, false\}$ extended with a new element $\perp$ for the undefined.

We will interpret each type $\sigma$ as a set $D(\sigma)$, and we let $D(ι) = \mathbb{N}_\perp$ and $D(o) = \mathbb{B}_\perp$.

Each of these sets will be partial orderings by letting $\perp$ be the smallest element and the rest maximal elements that have no ordering between them. Such ordered sets are called flat domains.

**Definition 4.5.11** A partial ordering $(D, \sqsubseteq)$ is called bounded complete if

1. Each bounded set has a least upper bound.

2. Each directed subset $X \subseteq D$ is bounded, with least upper bound $\sqcup X$. 

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It is easy to see that $\mathbb{N}_\bot$ and $\mathbb{B}_\bot$ both are bounded complete.

**Definition 4.5.12** Let $(D, \sqsubseteq_D)$ and $(E, \sqsubseteq_E)$ be two partial orderings that are bounded complete.

A function $f : D \to E$ is called *continuous* if $f$ is monotone and for all directed subsets $X$ of $D$ we have

$$\sqcup_D X = \sqcup_E \{ f(x) : x \in X \}.$$ 

**Lemma 4.5.13** Let $(D, \sqsubseteq)$ be bounded complete and let $f : D \to D$ be continuous. Then $f$ has a least fixed point in $D$.

*Proof*

The empty set has a least upper bound, which we call $\bot_D$.

Then $\bot_D \sqsubseteq_D f(\bot_D) \sqsubseteq_D f(f(\bot_D)) \sqsubseteq \cdots$.

The least upper bound of this sequence will be the least fixed point of $f$.

**Definition 4.5.14** Let $(D, \sqsubseteq_D)$ and $(E, \sqsubseteq_E)$ be two orderings that are bounded complete.

Let $D \to E$ be the set of continuous maps from $D$ to $E$.

If $f \in D \to E$ and $g \in D \to E$, we let $f \sqsubseteq_D \to E g$ if

$$\forall x \in D (f(x) \sqsubseteq g(x)).$$

**Lemma 4.5.15** Let $(D \to E, \sqsubseteq_D \to_E)$ be as above. This partial ordering will be bounded complete.

Not mentioning the ordering, we now interpret each type $\sigma$ by recursion on $\sigma$ as follows:

$$D(\sigma \to \tau) = D(\sigma) \to D(\tau).$$

What remains to be done can briefly be described as follows:

1. An *assignment* will be a map from a set of typed variables $x_\sigma^n$ to elements of $D(\sigma)$. The set of assignments can be viewed as elements of cartesian products of some $D(\sigma)$’s, and will be ordered in a bounded complete way by the coordinate-wise ordering.

2. Each term $M$ of type $\sigma$ will be interpreted as a continuous function $[M]$ from the set of assignments to $D(\sigma)$. It will sometimes be convenient to consider assignments restricted to the free variables of $M$, sometimes convenient to accept dummy variables.

3. We must show that the least fixed point operator is a continuous map from $(D, \sqsubseteq_D)$ to $D$. This is used to define $[Y_\sigma]$.

4. We must show that application is a continuous map from $(D \to E) \times D$ to $E$. This is used to define $[MN]$ from $[M]$ and $[N]$. 


5. We must show that abstraction is a continuous map from $D \times E \to F$ to $D \to (E \to F)$. This is used to define $[[\lambda x. M]]$ from $[[M]]$.

6. We must give interpretations to the constants $k_n$, $S$, $P$, $Z$, $\top$, and $\bot$.
This can easily be done by the reader.

7. We must show that if $M \longrightarrow N$ then $[[M]] = [[N]]$. The only case that requires some work is the reduction

$$((\lambda x. M)N) \longrightarrow M_N^N$$

where we must show by induction on $M$ that

$$[[M]](s_{[[N]]}(s)) = [[M_N^N]](s).$$

The proof by induction is not very hard.

The denotational semantics for PCF that we have given is the one used originally by Scott when he formed the logic LCF that was turned into the programming language PCF by Plotkin. There are however alternative ways of interpreting PCF-terms reflecting interesting aspects of PCF. One consequence of the existence of a denotational semantics is that a closed term of type $\top$ cannot be interpreted as two different numbers. Consequently, though the use of $\longrightarrow$ is a non-deterministic process, there is no canonical way of performing the next step, we do not risk to obtain different values to the same term.

We end this section by stating the converse, due to Plotkin, without any hints of proof:

**Theorem 4.5.16** Let $M$ be a closed term of type $\top$. If $[[M]] = n$, then $M \longrightarrow k_n$.

### 4.6 Exercises to Chapter 4

**Exercise 4.1** A **continued fraction** is a finite or infinite tree of fractions

$$\frac{1}{1 + n_0 + \frac{1}{1 + n_1 + \frac{1}{1 + n_2 + \cdots}}}.$$

A finite continued fraction will be a rational number while the value of an infinite continued fraction is defined as the limit of the finite subfractions. This limit will always exist. Why? (This part is not an exercise in logic, rather in elementary analysis.)

a) Show that if $0 < a < 1$, then there is a unique continued fraction with value $a$.

b) Show that the continued fraction of $a$ will be infinite if and only if $a$ is rational.
c) Show that the bijection obtained between the irrational elements in $[0, 1]$ and $\mathbb{N}^\mathbb{N}$ using continued fractions is a homeomorphism.

d) Discuss in which sense we may claim that the homeomorphism in c) is computable.

**Exercise 4.2** Show that if $F$ is a computable functional of type 2, then $F$ will have a computable associate.

**Exercise 4.3** Prove lemma 4.1.8.

**Exercise 4.4** Let $F : \mathbb{N}^\mathbb{N} \to \mathbb{N}$ be continuous, and let $T(F, f)$ be the set of finite sequences $\sigma$ bounded by $f$ such that $F$ is not constant on $B_\sigma$. Show that $T(F, f)$ will be a tree of sequences, and use König’s lemma to show that $T(F, f)$ is finite for each $f \in \mathbb{N}^\mathbb{N}$. Use this to give an alternative proof of Lemma 4.1.9.

**Exercise 4.5** Show that there is a continuous function $\hat{\Phi} : \mathbb{N}^\mathbb{N} \to \mathbb{N}^\mathbb{N}$ such that whenever $\alpha$ is an associate for a continuous $F : \mathbb{N}^\mathbb{N} \to \mathbb{N}$, then $\hat{\Phi}(\alpha)$ is an associate for $\Phi(F)$.

Show that $\hat{\Phi}$ can be chosen to be computable in the sense that $G(x, \alpha) = \hat{\Phi}(\alpha)(x)$ is computable.

**Exercise 4.6** Tait showed that the fan functional is not Kleene-computable. Fill out the details in the following proof of Tait’s result:

1. A quasi-associate for $F$ is a function $\alpha : \mathbb{N} \to \mathbb{N}$ such that
   - $\alpha(n) = m + 1 \Rightarrow F(f) = m$ when $\sigma_n \prec f$.
   - $\alpha(n) > 0 \land \sigma_n \prec \sigma_m \Rightarrow \alpha(m) = \alpha(n)$.
   - If $f \in 1 - sc(F)$ then for some $n$, $\sigma_n \prec f$ and $\alpha(n) > 0$.

If $\phi_e(\bar{a}, F) \downarrow$ and $\alpha$ is a quasi-associate for $F$, then there is some $n$ such that whenever $G$ has an associate extending $(\alpha(0), \ldots, \alpha(n))$ and $\phi_e(\bar{a}, G) \downarrow$ then

$$\phi_e(\bar{a}, G) = \phi_e(\bar{a}, F).$$

Hint: Use induction on the ordinal rank of the computation tree for $\phi_e(\bar{a}, F)$.

2. There is a quasi-associate for the constant zero function $^2O$ and a non-computable $f : \mathbb{N} \to \{0, 1\}$ such that $\alpha(f(n)) = 0$ for all $n$.

3. For any finite part $\bar{a}(n)$ of $\alpha$ there is an associate $\beta$ extending $\bar{a}(n)$ for a functional that is not constant on $\{0, 1\}^\mathbb{N}$.  

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4. Combining 2. and 3. we obtain a contradiction from the assumption that the fan functional is computable.

**Exercise 4.7** If $F : \mathbb{N}^\mathbb{N} \to \mathbb{N}$, let $F_i(f) = F((i, f))$.
Show that there is a total functional $\Gamma$ satisfying the equation

$$\Gamma(F) = F_0(\lambda x. \Gamma(F_{x+1})),$$

and that we may compute $\Gamma(F)$ uniformly from an associate for $F$.
The historical interest is that $\Gamma$ is not computable in the fan functional. This is too hard to be considered as an exercise.

**Exercise 4.8** Work out detailed proofs of Lemmas 4.2.5, 4.2.6 and 4.2.9.
This cannot be seen as a simple exercise, but rather as a minor project.

**Exercise 4.9** We reduce the definition of $\phi_e(\vec{a}, \vec{f}, \vec{F})$ by ignoring schemes v) and vi), and technically work with a subsystem.

a) Prove the $S_m^n$-theorem for computations with function and functional arguments as well as number arguments.

b) Prove the recursion theorem for computations with function and functional arguments as well as number arguments.

c) Use the recursion theorem to show that scheme v) is redundant in the original definition, i.e. that the class of functional computable in the subsystem is closed under primitive recursion.

d) Use the recursion theorem to show that scheme vi) is redundant in the original definition, i.e. that the $\mu$-operator is definable in the subsystem.

Hint: You may look at how the $\mu$-operator is defined in $PCF$.

**Exercise 4.10** Prove Lemma 4.4.10 and Lemma 4.4.11 in detail.

**Exercise 4.11** We actually have a hierarchy for sets definable by second order formulas as well. Let

- $A \subseteq X$ is $\Pi_{k+1}^1$ if for some $\Sigma_k^1$ set $B \subseteq \mathbb{N}^\mathbb{N} \times X$,
  $$\vec{x} \in A \Leftrightarrow \forall f \in \mathbb{N}^\mathbb{N}((f, \vec{x}) \in B).$$

- $A \subseteq X$ is $\Sigma_{k+1}^1$ if for some $\Pi_k^1$ set $B \subseteq \mathbb{N}^\mathbb{N} \times X$,
  $$\vec{x} \in A \Leftrightarrow \exists f \in \mathbb{N}^\mathbb{N}((f, \vec{x}) \in B).$$

a) Prove that the class of $\Pi_k^1$ sets are closed under finite unions and intersections, number quantifiers and universal function quantifiers. Prove that there are universal $\Pi_k^1$ sets for any $k \geq 1$ and any dimension.
b) Formulate and prove the analogue results for the \( \Sigma^1_k \)-classes.

**Exercise 4.12**  

a) Let \( f : \mathbb{N} \to \mathbb{N} \) be computable in \( 2^E \). Show that the image of \( f \) will also be computable in \( 2^E \). Discuss what this means for the relation between semicomputable sets and c.e. sets relative to \( 2^E \).

b) Show that the class of sets semicomputable in \( 2^E \) is closed under finite unions and intersections, but not under complements.

c) Show that two disjoint \( \Sigma^1_1 \)-sets can be separated by a \( \Delta^1_1 \)-set.  
   Hint: Use Gandy selection.

d) Show that there are disjoint \( \Pi^1_1 \)-sets that cannot be separated by any \( \Delta^1_1 \)-sets.  
   Hint: Use an analogue of the construction of two computably inseparable c.e. sets.

**Exercise 4.13**  

Let \( \sigma = \tau_1, \ldots, \tau_n \to b \) be a type, where \( b \) is a base type. Show that there is a term \( \supset \) of type \( o, \sigma, \sigma \to \sigma \) such that whenever \( s \) and \( t \) are terms of type \( \sigma \), then,

\[
\supset (tt)st \longrightarrow \lambda x_1 \tau_1 \cdot \cdots \cdot \lambda x_n \tau_n . sx_1 \cdots x_n.
\]

\[
\supset (ff)st \longrightarrow \lambda x_1 \tau_1 \cdot \cdots \cdot \lambda x_n \tau_n . tx_1 \cdots x_n.
\]

**Exercise 4.14** Prove Theorem 4.5.9 in detail.
Chapter 5

Non-trivial exercises and minor projects

In this chapter we have collected some exercises or minor projects that are based either on the material in more than one chapter or that are too hard to be considered as exercises in the ordinary sense. We call them exercises, though some of the items in this chapter are much more demanding than ordinary exercises.

**Exercise 5.1** Let $L$ be a finite language. Show that the relation ‘$\phi$ is true in all finite $L$-structures’ is decidable if all predicates are unary, while it in general is complete co-c.e.

**Exercise 5.2** Show that there is a computable partial ordering $(\mathbb{N}, \prec)$ such that every countable partial ordering $(X', <')$ can be embedded into $(\mathbb{N}, \prec)$.

**Exercise 5.3** In this exercise we will explore Skolem’s construction of a non-standard model for number theory. We adjust the proof to the setting that we have been using in this compendium.

Let $D^k$ be the set of subsets of $\mathbb{N}^k$ definable in the language of Theorem 1.2.2. We let $D = D^1$.

If $f : \mathbb{N} \to \mathbb{N}$, we let $f \in D^{func}$ if the graph of $f$ is in $D^2$.

a) Show that $D$ is countable.

b) Show that there is a maximal, nonprincipal filter $\mathcal{F}$ on $D$, i.e. a set $\mathcal{F} \subseteq D$ such that

- $N \in \mathcal{F}$.
- $\mathcal{F}$ has no finite elements.
- If $D \in \mathcal{F}$ and $D \subseteq E \in D$, then $E \in \mathcal{F}$.
- If $D \in \mathcal{F}$ and $E \in \mathcal{F}$, then $D \cap E \in \mathcal{F}$.
c) For \( f, g \in \mathcal{D}^{\text{func}} \), let
\[
f \equiv_{\mathcal{F}} g
\]
if
\[
\{ n : f(n) = g(n) \} \in \mathcal{F}.
\]
Show that \( \equiv_{\mathcal{F}} \) is an equivalence relation.

d) Let \( \mathcal{N}^{\ast} \) be the set of \( \equiv_{\mathcal{F}} \)-equivalence classes \([f]\), with the interpretations of the non-logical symbols in \( L \) as follows:

0\(^{\ast}\) and 1\(^{\ast}\) are the equivalence classes of the respective constant functions.

Addition, multiplication and exponentiation of equivalence classes are defined via the corresponding pointwise operations on representatives.

\([f] <^{\ast} [g] \iff \{ n : f(n) < g(n) \} \in \mathcal{F}\).

Work out the details and show that these concepts are well defined, i.e. independent of the choice of representatives.

e) Let \( \phi(x_1, \ldots, x_k) \) be a formula in \( L \).
Let \( f_1, \ldots, f_k \in \mathcal{D}^{\text{func}} \).
Show that
\[
\mathcal{N}^{\ast} \models \phi([f_1], \ldots, [f_k]) \iff \{ n : \mathcal{N} \models \phi(f_1(n), \ldots, f_k(n)) \} \in \mathcal{F}.
\]
Hint: Use induction on \( \phi \).

f) Show that \( \mathcal{N} \) and \( \mathcal{N}^{\ast} \) are not isomorphic, and describe an elementary embedding of \( \mathcal{N} \) into \( \mathcal{N}^{\ast} \).

**Exercise 5.4** For the sake of notational simplicity we extend the second order languages with variables for functions as well.

a) Find a second order formula that expresses that if \( X \) is a set, \( f : X \to X \) and \( x \in X \) then \( Y = \{ f^n(x) : n \in \mathbb{N} \} \).

b) Discuss to what extent our formal definition of a finite set is in accordance with your intuition about set theory.

c) It is a fact of set theory that every set accepts a total ordering. Use this and our formal definition of infinity to show that the following are equivalent.

1. \( X \) is finite.
2. Every total ordering of \( X \) is a well ordering.
Exercise 5.5 Let $L_=$ be the language of equality. We say that a subset $A$ of $\mathbb{N}$ is \textit{definable in second order finitary logic} if there is a sentence $\phi$ in $L^2_=$ such that $n \in A$ if and only if $\{0, \ldots, n - 1\} \models \phi$.

Discuss the properties of the class of sets definable in second order finitary logic in terms of decidability, complexity and closure properties.

Exercise 5.6 Let $\Sigma$ be a fixed alphabet. Show that there is an enumeration $\{M_i\}_{i \in \mathbb{N}}$ of the set of Turing machines over $\Sigma$ and a one-to-one enumeration $\{w_j\}_{j \in \mathbb{N}}$ of the set of words in $\Sigma^*$ such that the function $f$ satisfying the equation $f(i,j) \simeq k \iff M_i(w_j) = w_k$ is computable, where $\simeq$ means that either is both sides undefined, or both sides are defined and with the same value.

Exercise 5.7 In the text we have indicated how to establish some of the properties of the hierarchy $\{F_\alpha\}_{\alpha < \epsilon_0}$. Work out all the details and write a small essay about it. You may well combine this with Exercise 5.8.

Exercise 5.8 Show that if $\alpha < \epsilon_0$ then $F_\alpha$ is provably computable. Suggest a fundamental sequence for $\epsilon_0$ and thus a function $F_{\epsilon_0}$. If properly defined, $F_{\epsilon_0}$ will not be provably computable in PA, but in some second order extension. Discuss how we may formulate a proof of the totality of $F_{\epsilon_0}$ in some second order extension of PA.

Exercise 5.9 Let $L$ be a first order language, and let $X$ be a new predicate symbol of arity 1. Let $L'$ be $L$ extended with $X$. Let $\phi(x,X)$ be a formula in $L'$ with at most $x$ free.

The predicate $X$ (which we may think of as a set variable) is

\begin{itemize}
    \item \textit{positive in $\phi$} if $X$ does not occur in $\phi$, if $\phi$ is of the form $X(t)$, if $\phi$ is of one of the forms $\psi_1 \lor \psi_2$ or $\exists y \psi_1$ where $X$ is positive in $\psi_1$ and $\psi_2$ or if $\phi$ is of the form $\neg \psi$ where $X$ is negative in $\psi$.
    \item \textit{negative in $\phi$} if $X$ does not occur in $\phi$, if $\phi$ is of the form $\neg X(t)$, if $\phi$ is of one of the forms $\psi_1 \lor \psi_2$ or $\exists y \psi_1$ where $X$ is negative in $\psi_1$ and $\psi_2$ or if $\phi$ is of the form $\neg \psi$ where $X$ is positive in $\psi$.
\end{itemize}

Now, let $X$ be positive in $\phi(x,X)$. Let $\mathfrak{A}$ be an $L$-structure, $B \subseteq A$. Let $\Gamma_\phi(B) = \{a \in A : \phi(a,B)\}$.

a) Show that $\Gamma_\phi$ is a monotone operator on the set of subsets of $A$, i.e. that if $B \subseteq C$ then $\Gamma_\phi(B) \subseteq \Gamma_\phi(C)$.

b) Show that $\Gamma_\phi$ will have a least fixed point, i.e. there is a least set $B$ such that $\Gamma_\phi(B) = B$.

c) Now let $L$ be the language of number theory and $\mathfrak{A}$ the standard model $\mathbb{N}$ of number theory (i.e. the natural numbers as a structure). Let $X$ be positive in $\phi(x,X)$. Show that the least fixed point of $\Gamma_\phi$ is a $\Pi^1_1$-set.
d) Show that there is an $m$-maximal $\Pi^1_1$-set that can be defined by positive induction.

e) Hard! Show that there is a $\Pi^1_1$-set that is not the least fixed point of any $\Gamma_\phi$.

f) Optional! Show that if $X$ is positive in $\phi$, then $\Gamma_\phi$ will have a maximal fixed point.

Show that if $\phi$ is a positive $\Pi^0_1$-formula, i.e. we do positive induction over $\mathbb{N}$, then the maximal fixed point of $\Gamma_\phi$ will also be $\Pi^0_1$.

**Exercise 5.10** Let $L$ be a finite first order language with equality as the only relation symbol. Let $\mathfrak{A}$ be a countable $L$-structure.

We say that $\mathfrak{A}$ is computable if there is an onto map $\alpha : \mathbb{N} \to A$ and computable functions $\hat{f} : \mathbb{N}^k \to \mathbb{N}$ for each function symbol $f$ of arity $k$ such that

$$\alpha(\hat{f}(n_1, \ldots, n_k)) = F^\mathfrak{A}(\alpha(n_1), \ldots, \alpha(n_k))$$

for each function symbol $f$. 

$$\{(n, m) : \alpha(n) = \alpha(m)\}$$

is c.e.

**Warning!** Our definition is customized to this exercise, and the term ‘computable structure’ may be used in a different way elsewhere in the literature.

a) Trivial! Show that there exist computable fields of any characteristics.

b) Harder! Show that the algebraic closure of a computable field is itself computable.

c) Hard! Is there a computable field such that the set of irreducible polynomials is not computable? What is the best you can say about the complexity of the set of irreducible polynomials in a computable field?

**Exercise 5.11** Recall that a well ordering will be a total ordering $(X, <)$ such that any nonempty subset of $X$ will have a least element. An initial segment of a well ordering $(X, <)$ will be a subset $Y$ of $X$ such that $y \in Y \land x < y \Rightarrow x \in Y$, together with the ordering $< \text{ restricted to } Y$.

a) Show that an initial segment of a well ordering is itself a well ordering.

b) Show that if $(X_1, <_1)$ and $(X_2, <_2)$ are two well orderings and $\pi_1$ and $\pi_2$ are two isomorphisms between initial segments of $(X_1, <_1)$ and $(X_2, <_2)$ resp. , then $\pi_1(x) = \pi_2(x)$ if they are both defined.

Hint: Assume not. Consider the maximal initial segment where they agree, and show that they must agree on the ‘next’ element as well.

c) Show that if $(X_1, <_1)$ and $(X_2, <_2)$ are two well orderings then they are either isomorphic or one is isomorphic to a proper subset of the other. In all cases the isomorphism is unique.

Hint: Let $\pi$ be the union of all isomorphisms between initial segments of the two well orderings. Then $\pi$ is the unique isomorphism in question.
An ordinal number will be a set $\alpha$ such that

- $\alpha$ is transitive, i.e. if $\beta \in \alpha$ then $\beta$ is a set and if moreover $\gamma \in \beta$, then $\gamma \in \alpha$.
- $(\alpha, \in)$ is a well ordering.

(d) Show that the empty set $\emptyset$ is an ordinal number, and that $\emptyset \in \alpha$ for all non-empty ordinal numbers. Describe the three smallest ordinal numbers.

e) Show that if $\alpha$ is an ordinal number and $\beta \in \alpha$, then $\beta$ is an ordinal number.

f) Show that if $\alpha$ and $\beta$ are two ordinal numbers, then $\alpha = \beta$, $\alpha \in \beta$ or $\beta \in \alpha$.

Hint: Show that identity functions are the only possible isomorphisms between initial segments of ordinal numbers.

g) Show that if $\alpha$ is an ordinal number, then $\alpha + 1$ defined as $\alpha \cup \{\alpha\}$ is an ordinal number.

h) Show that the union of any set of ordinal numbers is an ordinal number.

A binary relation $(X, R)$ is well founded if each nonempty subset $Y$ of $X$ has an $R$-minimal element $x$, i.e. an element $x \in Y$ such that $y \not\in Y$ whenever $yRx$.

i) Let $(X, R)$ be a well founded relation. Let $R^*$ be the transitive (but not reflexive) closure of $R$. Show that $R^*$ is a well founded relation.

j) Show that if $(X, R)$ is a well founded relation, there is a unique function $\rho$ mapping $X$ into the ordinal numbers such that

$$\forall x \in X (\rho(x) = \bigcup \{\rho(y) + 1 ; yR^* x\}).$$

Hint: Use well foundedness to show that there is a maximal partial solution to the function equation for $\rho$, and once again to show that this maximal solution will be defined on all of $X$.

The function $\rho$ will be called the rank function of $R$. The image of $\rho$ will be an ordinal number and will be called the rank of $R$.

**Exercise 5.12** This exercise requires the combined knowledge of some axiomatic set theory, standard computability theory and computability in $2^E$.

Let $X$ be a hereditarily countable set, $A$ a binary relation on $B \subseteq \mathbb{N}$. We say that $A$ is a code for $X$ if $B, A$ is isomorphic to the transitive closure of $\{X\} \cup X$ with the $\in$-relation.

a) Show that the set of codes is $\Pi_1^1$.  

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b) Show that equality between the corresponding sets is a property of the codes computable in \( 2^E \). (Given two codes, \( 2^E \) can decide if they are codes for the same element or not.)

Hint: Use the recursion theorem for \( 2^E \).

c) Show that there are operators on the codes that are computable in \( 2^E \) and that simulates finite unions, unordered pairs and set subtraction on sets.

An ordinal number \( \alpha \) is called computable if \( \alpha \) has a computable code. Observe that if \( \alpha \) is computable and \( \beta < \alpha \), then \( \beta \) is computable.

d) Show that if \( \alpha \) is a computable ordinal number, then every \( X \in L(\alpha) \) will have a hyperarithmetical code.

(Here \( L(\alpha) \) is level no. \( \alpha \) in Gödel’s hierarchy of constructible objects.)

Let \( \omega^C_{1} \) (omega-1-Church-Kleene) be the least non-computable ordinal.

e) Show that the ordinal rank of a computable well founded tree will be computable.

Hint: Show that the partial ordering \( \sigma \leq \tau \Rightarrow \tau \prec \sigma \) on a computable well founded tree can be extended to a total, computable well ordering by combining it with the lexicographical ordering. (This combined ordering, if constructed correctly, is called the Kleene-Brouwer ordering).

f) Show that the ordinal rank of any \( \Delta^1_1 \)-well founded tree will be a computable ordinal.

Hint: Show that if this is not the case, we may use stage comparison to show that the set of computable well founded trees is computable in \( 2^E \). This will contradict that this set is uniformly \( \Pi^1_1 \) and that there are \( \Pi^1_1 \)-sets that are not hyperarithmetical.

g) Show that every hyperarithmetical set is in \( L(\omega^C_{1}) \).

Hint: Show that the set of computations of length \(< \alpha \), where the arguments may be \( 2^E \) and numbers, will be in \( L(\alpha + 1) \).

h) Show that \( L(\omega^C_{1}) \) satisfies \( \Delta_0 \)-comprehension and replacement.

Hint: For the replacement axiom you may find Gandy selection useful.

If you have worked through all items in this rather extensive exercise, you have proved that \( L(\omega^C_{1}) \) is a model of what is known as Kripke-Platek set theory, a much studied fragment of set theory.

**Exercise 5.13** Work out all the details in Section 4.5.3

**Exercise 5.14** Recall that a tree is a set of finite sequences, and that a tree will be well founded if there is no infinite branch.

A proof tree in an \( \omega \)-logic \( T \) will be a well founded tree \( T \) of sequences \( (\phi_0, \ldots, \phi_n) \) of formulas \( \phi_i \) with \( n \geq 0 \) such that
• If \((\phi_0, \ldots, \phi_n)\) and \((\phi'_0, \ldots, \phi'_m)\) are in \(T\), then \(\phi_0 = \phi'_0\).

• If \((\phi_0, \ldots, \phi_n)\) is a leaf node, then \(\phi_n\) is an axiom in \(T\).

• If \((\phi_0, \ldots, \phi_n)\) is not a leaf node, then \(\phi_n\) is a consequence of \(\{\phi ; (\phi_0, \ldots, \phi_n, \phi) \in T\}\) by the rules of \(\omega\)-logic.

By a suitable Gödel numbering, we may view the proof trees as trees over the natural numbers.

If \(T\) is a proof tree, we call it an \(\omega\)-proof for \(\phi_0\).

a) Show that for any formula \(\phi\), \(T \vdash_\omega \phi\) if and only if there is an \(\omega\)-proof for \(\phi_0\).

b) One technical obstacle in \(\omega\)-logic is that an \(\omega\)-proof may be infinite and consume all variables. Show that this is a minor obstacle, i.e. show that whenever \(y_1, \ldots, y_n\) are substitutable for \(x_1, \ldots, x_n\) in \(\phi\), and \(T \vdash_\omega \phi\), then \(T \vdash_\omega \phi_{y_1}^{x_1} \ldots y_n^{x_n}\).

Hint: Use induction on the ordinal rank of the proof tree.

Show that the proof tree of \(T_{y_1}^{x_1} \ldots y_n^{x_n}\) is computable in \(2^E\), the proof tree for \(\phi\) and the set of Gödel numbers for the axioms in \(T\).

c) Prove the theorem of constants for \(\omega\)-logic in detail.

Exercise 5.15 In the sections on model theory we focused on countable models and on finite models. This is of course only a part of model theory. In this exercise we will consider a tiny fragment of model theory for uncountable structures.

Definition 5.1 Let \(L\) be a first order language and \(\kappa\) a cardinal number. We say that an \(L\)-structure \(\mathfrak{A}\) is \(\kappa\)-saturated if for all sets \(B \subseteq A = |\mathfrak{A}|\) of cardinality \(< \kappa\) and all 1-types \(X\) in \(L[B]\) (relative to the theory of \(L[B]\)-sentences true in \(\mathfrak{A}\)) are realized in \(\mathfrak{A}\). \(\mathfrak{A}\) is fully saturated if \(\mathfrak{A}\) is \(\kappa\)-saturated where \(\kappa\) is the cardinality of \(|\mathfrak{A}|\).

Our definition of saturated corresponds to \(\omega\)-saturated.

a) Show that for any consistent theory \(T\) on \(L\) and cardinal number \(\kappa\) there exists a \(\kappa\)-saturated model for \(T\).

b) Assume the existence of an unbounded class of inaccessible cardinals.

Show that for any consistent theory \(T\) there exists a fully saturated model for \(T\) of arbitrary large cardinality.
Chapter 6

Appendix:
Some propositions from a beginners course in logic

Proposition 6.1 The Completeness Theorem, two versions

a) Let $T$ be a first order theory. Then $T$ is consistent if and only if $T$ has a model.

b) Let $T$ be a first order theory over the language $L$ and let $\phi$ be a formula in $L$. Then

$$T \vdash \phi \iff T \models \phi$$

Proposition 6.2 The Compactness Theorem

Let $T$ be a first order theory. Then $T$ has a model if and only if each finite subtheory $T_0$ of $T$ has a model.

Proposition 6.3 The Deduction Theorem

Let $T$ be a first order theory over a language $L$, let $\phi$ be a sentence in $L$ and $\psi$ a formula in $L$. Then

$$T, \phi \vdash \psi \iff T \vdash \phi \rightarrow \psi.$$ 

Proposition 6.4 The Theorem of Constants

Let $T$ be a first order theory over a language $L$ and let $c_1, \ldots, c_n$ be constants that are not in $L$. Let $L'$ be the extended language.

Let $T'$ be the theory over $L'$ with the same non-logical axioms as $T$, and let $\phi$ be a formula in $L$. Then for all lists $x_1, \ldots, x_n$ of variables

$$T \vdash \phi \iff T' \vdash \phi_{c_1, \ldots, c_n}.$$
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