# Categorial Grammars Based on Variants of the Lambek Calculus 

Candidate of Science (Ph. D.) Thesis

(translated from Russian)

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Moscow, 2012

Area: 01.01.06 - Mathematical Logic, Algebra, and Number Theory.
Defended on June 8, 2012 at Moscow State University (dissertation council No. D 501.001.84).
Confirmed by the Ministry of Education and Science of Russia on March 11, 2013 (order No. 130/nk).
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## Introduction

## Historical Background

The Lambek calculus L was introduced by J. Lambek [16] for defining natural language syntax using categorial grammars [22][21]. This calculus uses syntactic types built from primitive ones using three binary connectives: multiplication, left and right division.

Chomsky [7] suggested another family of grammars, called the Chomsky hierarchy. The most well-known class grammars from this hierarchy is the class of context-free grammars. Context-free grammars are widely used for parsing artificial languages (e.g., programming languages [1]), but for natural languages categorial grammars have significant advantages. In particular, they enjoy the lexicalisation property: all the information about syntax is kept in the categorial dictionary, and the analyser needs only that part of the dictionary, which is relevant to the text being parsed. Categorial grammars also allow to do semantic analysis, e.g., using Montague semantics [6].

We focus on comparing classes of languages generated by different classes of grammars just as sets of words, without syntactic structure. In this (so called weak) sense Lambek grammars are equivalent to context-free ones: the class of languages generated by grammars, based on L, coincides with the class of all context-free languages without the empty word [30]. Similar questions can be posed for grammars based on variants of the Lambek calculus (its fragments and extensions). It is known [5] that all context-free languages without the empty word can be generated by grammars based on the fragment of the Lambek calculus with only one division. Kanazawa [15] considered languages generated by grammars based on the Lambek calculus with additive conjunction and disjunction. This class of languages strictly contains finite intersections of context-free languages and lies inside the class of all context-sensitive languages. Moortgat [20] enriched the Lambek calculus with two modalities; as shown by Jäger [14], this extension does not enlarge the class of languages. Dikovsky and Dekhtyar [8] considered categorial dependency grammars (CDGs) based on a fragment of the Lambek calculus without multiplication, but with additional connectives for nonlocal dependencies, with a restriction: all the denominators are primitive types. Languages generated by CDGs form a special class that strictly contains the class context-free languages and is closed under finite unions, intersections with regular languages, and taking image or preimage of a homomorphism. As shown by Buszkowski [4], any recursively enumerable language can be generated by a grammar based on an extension of the Lambek calculus with a finite set of axioms.

The Lambek calculus is complete with respect to models interpreting Lambek types as formal languages over an alphabet (connectives of L correspond to multiplication, left and right divisions of languages) [31]. Such models are called L-models.

Derivability problems for L and its fragments $\mathrm{L}(\backslash, /)$ and $\mathrm{L}(\cdot, \backslash)$ are NP-complete as shown by Pentus [25] for L and by Savateev [32] for $\mathrm{L}(\backslash, /)$ and $\mathrm{L}(\cdot, \backslash)$. On the other hand, if the complexity of types is bounded (this is the usual case in practice), then there
are polynomial $\left(O\left(n^{5}\right)\right)$ time algorithms. These algorithms were independently presented by Pentus [26] and Fowler [9]. Fowler used such algorithms for parsing English sentences from CCGBank [10]. Derivability problem for $\mathrm{L}(\backslash)$ is decideable in polynomial time [32].

## Main Results

1. The class of languages generated by $\mathrm{L}_{\mathbf{1}}$-grammars coincides with the class of all context-free languages.
2. The derivability problem for $\mathrm{L}_{\mathbf{1}}(\backslash)$ is decidable in polynomial time.
3. The class of languages generated by $\mathrm{L}\left(\backslash ; p_{1}\right)$-grammars coincides with the class of all context-free languages without the empty word.
4. The derivability problem for $\mathrm{L}\left(\cdot, \backslash ; p_{1}\right)$ is NP-complete.
5. We present a calculus $\mathrm{L}^{\mathrm{R}}$ for the unary operation of language reversal and prove its completeness with respect to models on subsets of free semigroups (language models); the class of languages generated by $\mathrm{L}^{\mathrm{R}}$-grammars coincides with the class of all context-free languages without the empty word.

## Acknowledgments

The author thanks Prof. Mati Pentus for guiding him into the subject and constant attention to his work. He is also grateful to the members of the Department of Mathematical Logic and Theory of Algorithms, Moscow State University for the friendly atmosphere and constructive discussions.

## Chapter 1

## The Lambek Calculus and Categorial Grammars

### 1.1 The Calculus L and Its Fragments

We consider the Lambek calculus L, introduced in [16]. The countable set $\mathrm{Pr} \leftrightharpoons$ $\left\{p_{1}, p_{2}, p_{3}, \ldots\right\}$ is called the set of primitive types (here and further " $\leftrightharpoons$ " means "equals by definition"). Types of the Lambek calculus are built from primitive types using three binary connectives: <br>(left division), / (right division), and • (multiplication); we shall denote the set of all types by Tp. Formally Tp is defined as the smallest (by inclusion) set that satisfies the following two conditions:

1. $\operatorname{Pr} \subset \mathrm{Tp}$;
2. if $A, B \in \mathrm{Tp}$, then $(A \backslash B),(B / A),(A \cdot B) \in \mathrm{Tp}$.

Capital letters $(A, B, \ldots)$ range over types. Capital Greek letters range over finite (possibly empty) sequences of types; $\Lambda$ stands for the empty sequence; $A^{k}$ stands for a sequence of $A$ written $k$ times. Expressions of the form $\Gamma \rightarrow C$ are called sequents of the Lambek calculus; $\Gamma$ is the antecedent and $C$ succedent of $\Gamma \rightarrow C$. Sequents with empty antecedents are written as $\rightarrow C$.

The calculus L is defined by axioms $p_{i} \rightarrow p_{i}$ (called (ax)) for every $i$ and the following rules of inference:

$$
\begin{array}{ll}
\frac{A \Pi \rightarrow B}{\Pi \rightarrow A \backslash B}(\rightarrow \backslash), \text { where } \Pi \neq \Lambda ; & \frac{\Pi \rightarrow A \Gamma B \Delta \rightarrow C}{\Gamma \Pi(A \backslash B) \Delta \rightarrow C}(\backslash \rightarrow) ; \\
\frac{\Pi A \rightarrow B}{\Pi \rightarrow B / A}(\rightarrow /) \text {, where } \Pi \neq \Lambda ; & \frac{\Pi \rightarrow A \Gamma B \Delta \rightarrow C}{\Gamma(B / A) \Pi \Delta \rightarrow C}(/ \rightarrow) ; \\
\frac{\Gamma \rightarrow A \Delta \rightarrow B}{\Gamma \Delta \rightarrow A \cdot B}(\rightarrow \cdot) ; & \frac{\Gamma A B \Delta \rightarrow C}{\Gamma(A \cdot B) \Delta \rightarrow C}(\cdot \rightarrow) ; \\
\frac{\Pi \rightarrow A \Gamma A \Delta \rightarrow C}{\Gamma \Pi \Delta \rightarrow C}(\text { cut }) .
\end{array}
$$

Example 1.1. $\mathrm{L} \vdash\left(p_{1} \backslash p_{2}\right)\left(\left(p_{3} \backslash p_{2}\right) \backslash p_{4}\right) \rightarrow\left(\left(p_{3} \backslash p_{1}\right) \backslash p_{4}\right):$

$$
\begin{aligned}
& \frac{p_{3} \rightarrow p_{3} \quad \frac{p_{1} \rightarrow p_{1} \quad p_{2} \rightarrow p_{2}}{p_{1}\left(p_{1} \backslash p_{2}\right) \rightarrow p_{2}}}{\frac{p_{3}\left(p_{3} \backslash p_{1}\right)\left(p_{1} \backslash p_{2}\right) \rightarrow p_{2}}{\left(p_{3} \backslash p_{1}\right)\left(p_{1} \backslash p_{2}\right) \rightarrow\left(p_{3} \backslash p_{2}\right)}} \frac{\left(p_{3} \backslash p_{1}\right)\left(p_{1} \backslash p_{2}\right)\left(\left(p_{3} \backslash p_{2}\right) \backslash p_{4}\right) \rightarrow p_{4}}{\left(p_{1} \backslash p_{2}\right)\left(\left(p_{3} \backslash p_{2}\right) \backslash p_{4}\right) \rightarrow\left(\left(p_{3} \backslash p_{1}\right) \backslash p_{4}\right)}
\end{aligned}
$$

It is easy to see that sequents with empty antecedents cannot be derivable in L . The (cut) rule is eliminable:

Theorem 1 (J. Lambek, 1958). If a sequent is derivable in L , then it is derivable without using the (cut) rule. [16]

By $\operatorname{Tp}(\backslash)$ we denote the set of types not containing • and /. The calculus defined by axioms (ax) and rules $(\backslash \rightarrow)$ and $(\rightarrow \backslash)$ is called $\mathrm{L}(\backslash)$. (One can dually consider / instead of $\backslash$; we shall always use $\backslash$.) The calculi $\mathrm{L}(\cdot, \backslash)$ and $\mathrm{L}(\backslash, /)$ are defined in a similar way.

By $\operatorname{Tp}\left(p_{1}, \ldots, p_{N}\right)$ we denote the set of types containing only primitive types from the set $\left\{p_{1}, \ldots, p_{N}\right\}$. The corresponding fragment of the Lambek calculus is denoted by $\mathrm{L}\left(p_{1}, \ldots, p_{N}\right) ;$ its fragments with restricted sets of connectives are called $\mathrm{L}\left(\backslash ; p_{1}, \ldots, p_{N}\right)$, $\mathrm{L}\left(\cdot, \backslash ; p_{1}, \ldots, p_{N}\right), \mathrm{L}\left(\backslash, / ; p_{1}, \ldots, p_{N}\right)$. The most interesting case for us is $N=1$. For convenience we sometime write simply $p$ instead of $p_{1}$.

### 1.2 The Calculi $L^{*}$ and $L_{1}$

The calculus L* (the Lambek calculus allowing empty antecedents) is obtained from L by dropping the restriction $\Pi \neq \Lambda$ on the rules $(\rightarrow \backslash)$ and $(\rightarrow /)$.

By $\mathrm{Tp}_{1}$ we denote the set of types generated from primitive types and the constant 1 (unit) using the connectives $\backslash, /$, and $\cdot$. The calculus $\mathrm{L}_{\mathbf{1}}$ (the Lambek calculus with the unit, introduced in [17]) is obtained from L* by adding an extra axiom $\rightarrow \mathbf{1}$ (denoted by $(\rightarrow \mathbf{1}))$ and an extra rule

$$
\frac{\Gamma \Delta \rightarrow C}{\Gamma 1 \Delta \rightarrow C}(\mathbf{1} \rightarrow) .
$$

Fragments of $\mathrm{L}_{\mathbf{1}}$ and $\mathrm{L}^{*}$ with restricted sets of primitive types and/or connectives are defined exactly for $L$.

The cut elimination theorem [17] is true for $\mathrm{L}_{\mathbf{1}}$ and $\mathrm{L}^{*}$.

### 1.3 Conservativity. Substitution and Equivalence

Here and further $\mathscr{L}$ stands for one of the variants of the Lambek calculus: $\mathrm{L}, \mathrm{L}^{*}, \mathrm{~L}(\backslash)$, $\mathrm{L}^{*}(\backslash), \mathrm{L}\left(\backslash ; p_{1}, \ldots, p_{N}\right), \mathrm{L}^{*}\left(\backslash ; p_{1}, \ldots, p_{N}\right), \mathrm{L}_{\mathbf{1}}$, or $\mathrm{L}^{\mathrm{R}}$ (defined in Chapter 4). We denote the corresponding set of types ( Tp for L and $\mathrm{L}^{*}, \mathrm{Tp}(\backslash)$ for $\mathrm{L}(\backslash)$ etc) by $\mathrm{Tp}_{\mathscr{L}}$.

Definition. A calculus $\mathscr{L}_{1}$ is called a fragment of a calculus $\mathscr{L}_{2}$ if $\mathrm{Tp}_{\mathscr{L}_{1}} \subseteq \mathrm{Tp}_{\mathscr{L}_{2}}$ and any sequent built from types from $\operatorname{Tp}_{\mathscr{L}_{1}}$ is derivable in $\mathscr{L}_{1}$ if and only if it is derivable in $\mathscr{L}_{2}$. In this case $\mathscr{L}_{2}$ is called a conservative extension of $\mathscr{L}_{1}$.

Due to the cut elimination theorem $\mathrm{L}(\backslash)$ is a fragment of $\mathrm{L}, \mathrm{L}\left(\backslash ; p_{1}, \ldots, p_{N}\right)$ is a fragment of $\mathrm{L}(\backslash)$ etc. On the other hand, $\mathrm{L}^{*}$ is not a fragment of L . For example, the sequent $\left(p_{1} \backslash p_{1}\right) \backslash p_{2} \rightarrow p_{2}$ is derivable in $\mathrm{L}^{*}$, but not in L .

By $\mathscr{L}\left(p_{1}, \ldots, p_{N}\right)$ we denote the fragment of $\mathscr{L}$ with a restricted set of primitive types (e. g., L $\left.\left(\cdot, \backslash ; p_{1}, p_{2}, p_{3}\right)=\mathrm{L}(\cdot, \backslash)\left(p_{1}, p_{2}, p_{3}\right)\right)$.

The substitution of type $A$ for $z(z \in \operatorname{Pr} \cup\{\mathbf{1}\})$ into $\mathscr{A}$ is denoted by $\mathscr{A}[z:=A]$. Here $\mathscr{A}$ can be any syntactic object: a type, a sequence of types, a sequent, or (see further) a categorial grammar. The expression $\mathscr{A}\left[z_{1}:=A_{1}, z_{2}:=A_{2}, \ldots\right]\left(\right.$ or $\left.\mathscr{A}\left[z_{i}:=A_{i}\right]\right)$ means that all substitutions are performed simultaneously. In all calculi considered the substitution rule is admissible:

Proposition 1.1. If $\mathscr{L} \vdash \Pi \rightarrow C$, then $\mathscr{L} \vdash(\Pi \rightarrow C)\left[q_{1}:=A_{1}, q_{2}:=A_{2}, \ldots\right]$, where $q_{i} \in \operatorname{Pr}$.

Proof. All of the rules of our calculi are given by schemata, therefore they keep valid after substituting arbitrary types for primitive types (here it is essential, that $q_{i}$ is a primitive type, but not 1). It is sufficient to check that the axioms become derivable sequents, in other words, to derive $A \rightarrow A$ for every $A \in \mathrm{Tp}_{\mathscr{L}}$. This is done by induction on the length of type $A$. The base case is trivial, and induction step is established by the following derivations:

$$
\frac{A_{1} \rightarrow A_{1} A_{2} \rightarrow A_{2}}{\frac{A_{1} A_{2} \rightarrow A_{1} \cdot A_{2}}{A_{1} \cdot A_{2} \rightarrow A_{1} \cdot A_{2}}} \quad \frac{A_{1} \rightarrow A_{1} \quad A_{2} \rightarrow A_{2}}{\frac{A_{1}\left(A_{1} \backslash A_{2}\right) \rightarrow A_{2}}{A_{1} \backslash A_{2} \rightarrow A_{1} \backslash A_{2}}} \quad \frac{A_{1} \rightarrow A_{1} \quad A_{2} \rightarrow A_{2}}{\frac{\left(A_{2} / A_{1}\right) A_{1} \rightarrow A_{2}}{A_{2} / A_{1} \rightarrow A_{2} / A_{1}}}
$$

If for every $\Pi \rightarrow C$ the inverse implication also holds, the substitution is called faithful.

Definition. Types $A$ and $B$ are called equivalent in $\mathscr{L}(A \leftrightarrow \mathscr{L} B$ or just $A \leftrightarrow B)$ if $\mathscr{L} \vdash A \rightarrow B$ and $\mathscr{L} \vdash B \rightarrow A$.

The relation $\leftrightarrow \mathscr{L}$ is an equivalence relation. Moreover, it is a congruence w. r. t. the Lambek calculus connectives:

Proposition 1.2. The relation $\leftrightarrow \mathscr{L}$ is reflexive, symmetric and transitive. If $A_{1} \leftrightarrow A_{2}$ and $B_{1} \leftrightarrow B_{2}$, then $A_{1} \cdot B_{1} \leftrightarrow A_{2} \cdot B_{2}, A_{1} \backslash B_{1} \leftrightarrow A_{2} \backslash B_{2}, B_{1} / A_{1} \leftrightarrow B_{2} / A_{2}$.

Proof. Reflexivity follows from the derivability of $A \rightarrow A$, transitivity is due to (cut), symmetricity is obvious.

The remaining statements are proved using, respectively, the rules $(\cdot \rightarrow)$ and $(\rightarrow \cdot)$, $(\rightarrow \backslash)$ and $(\backslash \rightarrow),(\rightarrow /)$ and $(/ \rightarrow)$.

Using (cut) we prove that if replacing a subtype with an equivalent one does not affect derivability.

### 1.4 Lambek Categorial Grammars

We call an alphabet an arbitrary nonempty finite set. The set of all finite sequences (including the empty one) of elements of an alphabet $\Sigma$ (words over $\Sigma$ ) is denoted by $\Sigma^{*}$. The empty word is denoted by $\epsilon$. The set of all nonempty word is denoted by $\Sigma^{+}$. Subsets of $\Sigma^{*}$ are called formal languages (or simply languages) over $\Sigma$.

Languages are usually infinite as sets. Various types of formal grammars are used to define some of them in a finite way. The Lambek calculus and its variants are the base for Lambek categorial grammars.

Definition. A grammar based on the calculus $\mathscr{L}$ (an $\mathscr{L}$-grammar) is a triple $\mathcal{G}=$ $\langle\Sigma, H, \triangleright\rangle$, where $\Sigma$ is an alphabet, $H \in \mathrm{Tp}_{\mathscr{L}}$, and $\triangleright \subset \mathrm{Tp}_{\mathscr{L}} \times \Sigma$ is a finite binary correspondence. The language generated by $\mathcal{G}$ is $\mathfrak{L}(\mathcal{G}) \leftrightharpoons\left\{a_{1} \ldots a_{k} \in \Sigma^{*} \mid\right.$ $\left.\exists B_{1}, \ldots, B_{k}: B_{i} \triangleright a_{i} \mathscr{L} \vdash B_{1} \ldots B_{k} \rightarrow H\right\}$. Such languages are called $\mathscr{L}$-languages.

The notion of substitution into $\mathscr{L}$-grammar is defined in a natural way: if $\mathcal{G}=$ $\langle\Sigma, H, \triangleright\rangle$ is an $\mathscr{L}$-grammar, $z \in \operatorname{Pr} \cup\{\mathbf{1}\}$ and $A \in \operatorname{Tp}_{\mathscr{L}}$, then $\mathcal{G}[z:=A] \leftrightharpoons\langle\Sigma, H[z:=$ $A], \triangleright[z:=A]\rangle$, where $\triangleright[z:=A] \leftrightharpoons\{\langle B[z:=A], a\rangle \mid B \triangleright a\}$.
Proposition 1.3. If $\mathscr{L}$ is a fragment of $\mathscr{L}^{\prime}$, then any $\mathscr{L}$-language is an $\mathscr{L}^{\prime}$-language.
Proof. Any $\mathscr{L}$-grammar can be considered as $\mathscr{L}^{\prime}$-grammar. The language will remain the same due to conservativity.

Vice versa, if an $\mathscr{L}^{\prime}$-grammar contains only types from $\mathrm{Tp}_{\mathscr{L}}$, it can be considered as an $\mathscr{L}$-grammar generating the same language.

Besides categorial grammars, based on the Lambek calculus and its variants, languages can be also defined by a family of formalisms called the Chomsky hierarchy [7]. We consider Chomsky grammars of type 2, also called context-free grammars.

Definition. A context-free grammar is a quadruple $G=\langle N, \Sigma, P, S\rangle$, where $N$ and $\Sigma$ are two disjoint alphabets, $P \subset N \times(N \cup \Sigma)^{*}, P$ is finite, and $S \in N$. We define a binary relation $\Rightarrow_{G}$ as follows: for all $\omega, \psi \in(N \cup \Sigma)^{*}$ we have $\omega \Rightarrow_{G} \psi$ if and only if $\omega=\eta A \theta$, $\psi=\eta \beta \theta$, and $\langle A, \beta\rangle \in P$ for some $A \in N, \beta, \eta, \theta \in(N \cup \Sigma)^{*}$. The binary relation $\Rightarrow_{G}^{*}$ is the reflexive transitive closure of $\Rightarrow_{G}$. The language $\mathfrak{L}(G)=\left\{w \in \Sigma^{*} \mid S \Rightarrow_{G}^{*} w\right\}$ is the language generated by $G$. Such languages are called context-free.

We study the relations between classes of languages generated, on one hand, by context-free grammars, and, on the other hand, by $\mathscr{L}$-grammars for various $\mathscr{L}$. The answers for $\mathscr{L}=\mathrm{L}$ and $\mathscr{L}=\mathrm{L}(\backslash)$ are given by the following theorems:

Theorem 2 (C. Gaifman, 1960; W. Buszkowski, 1985). Every context-free language without the empty word is an $\mathrm{L}(\backslash)$-language. $[2,5]$

Theorem 3 (M. Pentus, 1993). Every L-language is context-free and does not contain the empty word. [30]

Proposition 1.3 now yields that the following three classes of languages coincide: L-languages, L( $\backslash$ )-languages, and context-free languages without the empty word.

These theorems have variants for the calculi allowing empty antecedents. Theorem 3 remains true:

Theorem 4 (M. Pentus, 1993). Every L*-language is context-free. [30]
To prove the inverse statement we shall need a refinement of Theorem 2. Essentially, it follows from the Greibach normal form theorem for context-free grammars.

Definition. Let us call a type a $G$-type (after Greibach and Gaifman) if it is of the form $p_{i}, p_{j} \backslash p_{i}$, or $p_{k} \backslash\left(p_{j} \backslash p_{i}\right)$, and $k, j \neq 1$. Let us call a sequent $\Pi \rightarrow C$ a $G$-sequent if all types in $\Pi$ are G-types and $C=p_{1}$.

Theorem 5. Any context-free language without the empty word is generated by some $\mathrm{L}(\backslash)$-grammar $\mathcal{G}=\langle\Sigma, H, \triangleright\rangle$, where $H=p_{1}$ and all types used in $\triangleright$ are $G$-types.

Evidently, when we check whether a word is in a language generated by such sort of grammar, all sequents we try to derive are G-sequents.

Proposition 1.4. If $\Pi \rightarrow p_{1}$ is a $G$-sequent, then

$$
\mathrm{L}^{*} \vdash \Pi \rightarrow p_{1} \Longleftrightarrow \mathrm{~L} \vdash \Pi \rightarrow p_{1} .
$$

Proof. The succedents of all sequents in a derivation of a G-sequent (both in L and in $\left.L^{*}\right)$ are primitive types. Therefore, the rules $(\rightarrow \backslash)$ and $(\rightarrow /)$ cannot be applied, and other rules and axioms are the same in L and $\mathrm{L}^{*}$.

Now let $M$ be a context-free language. If $\epsilon \notin M$, then let $\mathcal{G}$ be the $\mathrm{L}(\backslash)$-grammar for $M$ given by Theorem 5. Due to Proposition 1.4 the language $\mathfrak{L}(\mathcal{G})$ will not change if we consider $\mathcal{G}$ as an $L^{*}(\backslash)$-grammar. The case when $\epsilon \in M$ is handled in Section 1.8.

Definition. Categorial grammars in which $\triangleleft$ is a partial function from $\Sigma$ to $\operatorname{Tp}_{\mathscr{L}}$ (in other words, for every $a \in \Sigma$ there exists not more than one type $B \in \operatorname{Tp}_{\mathscr{L}}$, such that $B \triangleright a$ ), are called grammars with single type assignment. (The term"deterministic grammars", used in in [33], is not as appropriate, since these grammars still have some nondeterminicity: the sequent $B_{1} \ldots B_{n} \rightarrow H$ can have several essentially different derivations in $\mathscr{L}$.)

Theorem 6 (A. Safiullin, 2007). Every context-free language without the empty word can be generated by an L-grammar with single type assignment. [33]

In [33] it is also stated that this theorem remains true for $\mathrm{L}(\backslash, /)$. The question for other fragments of L (namely, for $\mathrm{L}(\backslash)$ or even for $\mathrm{L}(\cdot, \backslash)$ ) and for $\mathrm{L}_{\mathbf{1}}$ and its fragments remains open.

### 1.5 Complexity Results for the Lambek Calculus and its Fragments

In this section we give a brief survey of known results concerning complexity of derivation problems for the Lambek calculus and its fragments. A more substantial discussion can be found in [27].

Due to the cut elimination theorem every derivable sequent has a derivation of polynomial (w. r. t. the length of the sequent) size. Therefore, derivability problems for $\mathrm{L}, \mathrm{L}_{1}$ and their fragments are in the NP class.

Theorem 7 (M. Pentus, 2003). Derivability problems for L and L* (and, therefore, also for $\mathrm{L}_{1}$ ) are NP-complete. [25]

A stronger result was obtained by Savateev:
Theorem 8 (Yu. Savateev, 2008-2009). Derivability problems for $\mathrm{L}(\backslash, /)$, L* $(\backslash, /)$, $\mathrm{L}(\cdot, \backslash)$, and $\mathrm{L}^{*}(\cdot, \backslash)$ are $N P$-complete.

The case of only one connective is different:
Theorem 9 (Yu. Savateev, 2006). There exists an algorithm that checks derivability of a sequent in $\mathrm{L}(\backslash)$ ( $\mathrm{L}^{*}(\backslash)$ ) in $O\left(n^{3}\right)$ steps, where $n$ is the length of the sequent (the number of connective and primitive type occurrences).

Moreover, there eixsts (Yu. Savateev, 2008) an algorithm that checks whether a word belongs to the language generated by a given $\mathrm{L}(\backslash)$-grammar ( $\mathrm{L}^{*}(\backslash)$-grammar) in polynomial time: the number of steps is not greater than a polynom of the sum of the length of the word and the size of the grammar.

Savateev's results can be found in his C. Sc. (Ph. D.) thesis [32].

### 1.6 Multiplicative Cyclic Linear Logic

We consider an auxiliary calculus $\mathrm{MCLL}_{1, \perp}$, the multiplicative linear logic. $\mathrm{MCLL}_{1, \perp}$ is a fragment of Girard's [11] linear logic and was first mentioned in [28]. Its connection with the Lambek calculus is established in [23].

Elements of a countable set $\operatorname{Var}=\left\{p_{1}, p_{2}, p_{3}, \ldots\right\}$ are called variables; At $\leftrightharpoons \operatorname{Var} \cup$ $\{\bar{q} \mid q \in \operatorname{Var}\} \cup\{\mathbf{1}, \perp\}$ is the set of atoms. Formulae of $\mathrm{MCLL}_{\mathbf{1}, \perp}$ are built from atoms using two binary connectives: $>($ multiplicative disjunction, "par") and $\otimes$ (multiplicative conjunction, "tensor"). We denote the set of all formulae by $\mathrm{Fm}_{1, \perp}$. Capital Latin letters range over formulae. Capital Greek letters denote finite sequences of formulae; $\Lambda$ stands for the empty sequence. Sequents of $\mathrm{MCLL}_{1, \perp}$ are of the form $\rightarrow \Gamma$.

Definition. The linear negation is introduced externally as an inductively defined mapping $(\cdot)^{\perp}: \mathrm{Fm} \rightarrow \mathrm{Fm}: p_{i}^{\perp} \leftrightharpoons \bar{p}_{i}, \bar{p}_{i}^{\perp} \leftrightharpoons p_{i},(A \ngtr B)^{\perp} \leftrightharpoons B^{\perp} \otimes A^{\perp},(A \otimes B)^{\perp} \leftrightharpoons B^{\perp} 8 A^{\perp}$.

Axioms of $\mathrm{MCLL}_{1, \perp}$ are sequents of the form $\rightarrow p_{i} \bar{p}_{i}$ and the sequent $\rightarrow \mathbf{1}$. Rules of inference:

$$
\begin{aligned}
\frac{\rightarrow \Gamma A B \Delta}{\rightarrow \Gamma(A>B) \Delta}(\rightarrow 8) ; & \frac{\rightarrow \Gamma A \rightarrow B \Delta}{\rightarrow \Gamma(A \otimes B) \Delta}(\rightarrow \otimes) ; \\
\frac{\rightarrow \Gamma \Delta}{\rightarrow \Gamma \perp \Delta}(\rightarrow \perp) ; & \frac{\rightarrow \Gamma \Delta}{\rightarrow \Delta \Gamma}(\mathrm{rot}) .
\end{aligned}
$$

We also consider two-sided $\mathrm{MCLL}_{1, \perp}$-sequents: $A_{1} A_{2} \ldots A_{n} \rightarrow B_{1} \ldots B_{m}$ means $\rightarrow A_{n}^{\perp} \ldots A_{2}^{\perp} A_{1}^{\perp} B_{1} \ldots B_{m}$.

The fragment of $\mathrm{MCLL}_{1, \perp}$ without constants (defined by axioms $\rightarrow p_{i} \bar{p}_{i}$ and rules $(\rightarrow \varnothing),(\rightarrow \otimes)$, and (rot)) is called MCLL.

Definition. The standard translation $\widehat{A} \in \mathrm{Fm}$ of $A \in \mathrm{Tp}_{1}$ is defined as follows:

1. $\widehat{p}_{i} \leftrightharpoons p_{i}\left(p_{i} \in \operatorname{Pr}\right)$;
2. $\widehat{\mathbf{1}} \leftrightharpoons \mathbf{1}$;
3. $\widehat{A \cdot B} \leftrightharpoons \widehat{A} \otimes \widehat{B}$;
4. $\widehat{A \backslash B} \leftrightharpoons \widehat{A}^{\perp} \varnothing \widehat{B}$;
5. $\widehat{B / A} \leftrightharpoons \widehat{B} \gtrdot \widehat{A}^{\perp}$.

If $\Pi=B_{1} \ldots B_{m}$, then $\widehat{\Pi} \leftrightharpoons \widehat{B}_{1} \ldots \widehat{B}_{m}$. The standard translation of a sequent $\Pi \rightarrow C$ is the sequent $\widehat{\Pi} \rightarrow \widehat{C}$, in other notation, $\rightarrow \widehat{B}_{m}^{\perp} \ldots \widehat{B}_{1}^{\perp} C$.

In terms of this translation $\mathrm{MCLL}_{\mathbf{1}, \perp}$ is a conservative extension of $\mathrm{L}_{\mathbf{1}}$ :
Theorem 10. Let $\Pi=B_{1} \ldots B_{m}$ and $B_{1}, \ldots, B_{m}, C \in \mathrm{Tp}_{\mathbf{1}}$. Then

$$
\mathrm{L}_{1} \vdash \Pi \rightarrow C \Longleftrightarrow \mathrm{MCLL}_{1, \perp} \vdash \widehat{\Pi} \rightarrow \widehat{C}
$$

Due to conservativity of $\mathrm{MCLL}_{\mathbf{1}, \perp}$ over MCLL and $\mathrm{L}_{\mathbf{1}}$ over $\mathrm{L}^{*}$ we get a similar theorem for $L^{*}$ :

Theorem 11. Let $\Pi=B_{1} \ldots B_{m} \quad B_{1}, \ldots, B_{m}, C \in \mathrm{Tp}$. Then

$$
\mathrm{L}^{*} \vdash \Pi \rightarrow C \Longleftrightarrow \mathrm{MCLL} \vdash \widehat{\Pi} \rightarrow \widehat{C}
$$

The calculus L corresponds to a modified version of MCLL, called MCLL'. It is obtained from MCLL by adding the restriction $\Gamma \Delta \neq \Lambda$ on the $(\rightarrow \gamma)$ rule.

Theorem 12. Let $\Pi=B_{1} \ldots B_{m} \quad B_{1}, \ldots, B_{m}, C \in \mathrm{Tp}$. Then

$$
\mathrm{L} \vdash \Pi \rightarrow C \Longleftrightarrow \mathrm{MCLL}^{\prime} \vdash \widehat{\Pi} \rightarrow \widehat{C} .
$$

Theorems 10 and 12 are proved in [23].
Example 1.2. In this example we write $A^{\wedge}$ instead of $\widehat{A}$. By definition $\left(\left(p_{3} \backslash p_{1}\right) \backslash p_{4}\right)^{\wedge}=$ $\left(\bar{p}_{1} \otimes p_{3}\right) \ngtr p_{4},\left(\left(p_{3} \backslash p_{2}\right) \backslash p_{4}\right)^{\wedge \perp}=\bar{p}_{4} \otimes\left(\bar{p}_{3} \diamond p_{2}\right)$, and $\left(p_{1} \backslash p_{2}\right)^{\wedge \perp}=\bar{p}_{2} \otimes p_{1}$. Therefore, by Example 1.1 and Theorem 12 we get MCLL $\vdash \rightarrow\left(\bar{p}_{4} \otimes\left(\bar{p}_{3} \diamond p_{2}\right)\right)\left(\bar{p}_{2} \otimes p_{1}\right)\left(\left(\bar{p}_{1} \otimes p_{3}\right) \ngtr p_{4}\right)$. Here is the explicit derivation:

$$
\begin{aligned}
& \frac{\rightarrow p_{2} \bar{p}_{2} \rightarrow p_{1} \bar{p}_{1}}{} \frac{\rightarrow p_{2}\left(\bar{p}_{2} \otimes p_{1}\right) \bar{p}_{1} \rightarrow p_{3} \bar{p}_{3}}{\rightarrow p_{2}\left(\bar{p}_{2} \otimes p_{1}\right)\left(\bar{p}_{1} \otimes p_{3}\right) \bar{p}_{3}} \\
& \rightarrow \frac{p_{4} \bar{p}_{4}}{\rightarrow \bar{p}_{3} p_{2}\left(\bar{p}_{2} \otimes p_{1}\right)\left(\bar{p}_{1} \otimes p_{3}\right)} \\
& \frac{\rightarrow\left(\bar{p}_{3} \& p_{2}\right)\left(\bar{p}_{2} \otimes p_{1}\right)\left(\bar{p}_{1} \otimes p_{3}\right)}{\rightarrow p_{4}\left(\bar{p}_{4} \otimes\left(\bar{p}_{3} \& p_{2}\right)\right)\left(\bar{p}_{2} \otimes p_{1}\right)\left(\bar{p}_{1} \otimes p_{3}\right)} \\
& \rightarrow\left(\bar{p}_{4} \otimes\left(\bar{p}_{3} \oslash p_{2}\right)\right)\left(\bar{p}_{2} \otimes p_{1}\right)\left(\bar{p}_{1} \otimes p_{3}\right) p_{4} \\
& \rightarrow\left(\bar{p}_{4} \otimes\left(\bar{p}_{3} \& p_{2}\right)\right)\left(\bar{p}_{2} \otimes p_{1}\right)\left(\left(\bar{p}_{1} \otimes p_{3}\right) \ngtr p_{4}\right)
\end{aligned}
$$

Substitution $\rightarrow \Gamma[z:=A]$, where $z \in \operatorname{Var} \cup\{\mathbf{1}\}$, is defined, respecting the negation, in the following way: if $z=p_{i}$, then for every $p_{i}$ we substitute $A$ and for every $\bar{p}_{i}$ we substitute $A^{\perp}$. If $z=\mathbf{1}$, then $A$ is substituted for each occurrence of $\mathbf{1}$, and $A^{\perp}$ is substituted for each occurrence of $\perp$. The substitution rule is admissible in both $\mathrm{MCLL}_{1, \perp}$ and MCLL':

Proposition 1.5. If $\rightarrow \Gamma$ is derivable in $\mathrm{MCLL}_{1, \perp}$ (resp., in $\mathrm{MCLL}^{\prime}$ ), then $\rightarrow \Gamma\left[q_{1}:=\right.$ $\left.A_{1}, q_{2}:=A_{2}, \ldots\right]$, where $q_{i} \in \operatorname{Var}, A_{i} \in \mathrm{Fm}$, is also derivable in $\mathrm{MCLL}_{\mathbf{1}, \perp}$ (resp., MCLL').

The proof is similar as for the Lambek calculus (Proposition 1.1).

### 1.7 Proof Nets

We introduce proof nets (graph-theoretical criteria of derivability) for MCLL and MCLL'. The idea of proof nets belongs to mathematical folklore. We consider the variant of proof nets by Pentus [24] and its adaptation for MCLL'. This variant of proof nets is also close to the one proposed by de Groote [12]. Due to Theorems 11 and 12 these criteria can also be used to study derivability in $\mathrm{L}^{*}$ and L (and this is our goal).

Let $\rightarrow \Gamma$ be a sequent of MCLL or MCLL'. First we build a relation structure $\boldsymbol{\Omega}_{\Gamma}=\left\langle\Omega_{\Gamma},<_{\Gamma}, \prec_{\Gamma}\right\rangle$. Let $\Gamma=B_{1} \ldots B_{m}$. We put $\diamond$ signs before $B_{1}$ and between $B_{i}$ and $B_{i+1}(i=1, \ldots, m-1)(\diamond$ is a new formal symbol, $\diamond \notin \mathrm{Fm}): \diamond B_{1} \diamond B_{2} \diamond \cdots \diamond B_{m}$. In this string we number all symbols except brackets (an atom is considered one symbol) and denote the set of pairs 〈symbol, number〉 by $\Omega_{\Gamma}$. Elements of $\Omega_{\Gamma}$ are called occurrences of the corresponding symbols in $\Gamma$ and denoted by lowercase Greek letters. For $\alpha=\left\langle s_{1}, k_{1}\right\rangle$, $\beta=\left\langle s_{2}, k_{2}\right\rangle \in \Omega_{\Gamma}$ we define $\alpha<_{\Gamma} \beta \Longleftrightarrow k_{1}<k_{2}$.

For a subformula of $\Gamma$ the corresponding subset of $\Omega_{\Gamma}$ is called the occurrence of this subformula ( $B_{1}, B_{2}, \ldots, B_{m}$ are also considered subformulae). If $X$ is a subformula occurrence, then we denote by $\mathrm{l}(X)$ the occurrence of a connective or $\diamond$ immediately to the left of $X$ and by $\mathrm{r}(X)$ the occurrence of a connective or $\diamond$ immediately to the right
of $X$ (if there is no such occurrence, $\mathrm{r}(X)$ is defined cyclically as the leftmost occurrence of $\diamond)$.

The transitive closure of the union of parse trees of the formulae $B_{1}, \ldots, B_{m}$ (where the vertices are occurrences, that is, elements of $\Omega_{\Gamma}$ ) is denoted by $\prec_{\Gamma}$. In other words, $\alpha \prec_{\Gamma} \beta$ if and only if $\alpha$ and $\beta$ are the occurrences of main connectives of two formulae (or atoms), such that the first one is a subformula of the second one.

We denote the set of all occurrences of 8 by $\Omega_{\Gamma}^{\diamond}$, of $\otimes$ by $\Omega_{\Gamma}^{\otimes}$, of $\diamond$ by $\Omega_{\Gamma}^{\diamond}$, of $p_{1}, p_{2}, \ldots$ by $\Omega_{\Gamma}^{\mathrm{At}^{+}}$, of $\bar{p}_{1}, \bar{p}_{2}, \ldots$ by $\Omega_{\Gamma}^{\mathrm{At}^{-}} ; \Omega_{\Gamma}^{8 \diamond} \leftrightharpoons \Omega_{\Gamma}^{8} \cup \Omega_{\Gamma}^{\diamond}, \Omega_{\Gamma}^{\otimes \diamond \diamond} \leftrightharpoons \Omega_{\Gamma}^{\otimes} \cup \Omega_{\Gamma}^{8 \diamond}, \Omega_{\Gamma}^{\mathrm{At}^{8 t}} \leftrightharpoons \Omega_{\Gamma}^{\mathrm{At}^{+}} \cup \Omega_{\Gamma}^{\mathrm{At}^{-}}$. For $X \subseteq \Omega_{\Gamma}$ we define the count check $\#(X) \leftrightharpoons\left|X \cap \Omega_{\Gamma}^{\mathrm{At}^{+}}\right|-\left|X \cap \Omega_{\Gamma}^{\mathrm{At}^{-}}\right|$.

For $\alpha, \beta \in \Omega_{\Gamma}$ we define $(\alpha, \beta) \leftrightharpoons\left\{\gamma \mid \alpha<_{\Gamma} \gamma<_{\Gamma} \beta\right\}, \operatorname{In}(\alpha, \beta) \leftrightharpoons(\alpha, \beta) \cup(\beta, \alpha)$, and $\operatorname{Out}(\alpha, \beta) \leftrightharpoons \Omega_{\Gamma}-(\operatorname{In}(\alpha, \beta) \cup\{\alpha, \beta\})$ (the minus sign here means set-theoretic difference).

Definition. An oriented graph $\left\langle\Omega_{\Gamma}, \mathcal{C}\right\rangle$, where $\mathcal{C} \subseteq \Omega_{\Gamma} \times \Omega_{\Gamma}$, is called $<_{\Gamma}$-planar if for any $\langle\alpha, \beta\rangle \in \mathcal{C}$ and $\langle\gamma, \delta\rangle \in \mathcal{C}$ such that $\{\alpha, \beta\} \cap\{\gamma, \delta\}=\emptyset$ we have $\gamma \in \operatorname{In}(\alpha, \beta) \Longleftrightarrow \delta \in$ $\operatorname{In}(\alpha, \beta)$. Geometrically this means that the graph can be drawn in the upper semiplane without intersections if its vertices lie on the semiplane's border in the $<_{\Gamma}$ order.

Definition. A proof net is a triple $\mathfrak{N}=\left\langle\Omega_{\Gamma}, \mathcal{A}, \mathcal{E}\right\rangle$, where $\mathcal{A} \subset \Omega_{\Gamma} \times \Omega_{\Gamma}$ and $\mathcal{E} \subset \Omega_{\Gamma} \times \Omega_{\Gamma}$, that satisfies the following axioms:

1. $\left|\Omega_{\Gamma}^{8 \diamond}\right|-\left|\Omega_{\Gamma}^{\otimes}\right|=2$;
2. $\mathcal{A}$ is a total function from $\Omega_{\Gamma}^{\otimes}$ to $\Omega_{\Gamma}^{8 \diamond}$;
3. $\mathcal{E}$ is a bijective function from $\Omega_{\Gamma}^{\mathrm{A}^{+}}$to $\Omega_{\Gamma}^{\mathrm{At}^{-}}$, and if $\alpha$ is an occurrence of $p_{i}$ then $\mathcal{E}(\alpha)$ is an occurrence of $\bar{p}_{i}$;
4. the graph $\left\langle\Omega_{\Gamma}, \mathcal{A} \cup \mathcal{E}\right\rangle$ is $<_{\Gamma}$-planar;
5. the graph $\left\langle\Omega_{\Gamma}, \mathcal{A} \cup \prec_{\Gamma}\right\rangle$ is acyclic.

Definition. A strong proof net is a proof net that additionally satisfies the following axiom:
6. for any subformula occurrence $X \subset \Omega_{\Gamma}$ we have $\widetilde{\mathcal{A}}(1(X)) \neq \widetilde{\mathcal{A}}(\mathrm{r}(X))$, where the mapping $\widetilde{\mathcal{A}}: \Omega_{\Gamma}^{\otimes \otimes \diamond} \rightarrow \Omega_{\Gamma}^{8 \diamond}$ is defined as follows:

$$
\widetilde{\mathcal{A}}(\alpha) \leftrightharpoons \begin{cases}\alpha & \text { if } \alpha \in \Omega_{\Gamma}^{8 \diamond}, \\ \mathcal{A}(\alpha) & \text { if } \alpha \in \Omega_{\Gamma}^{\otimes}\end{cases}
$$

Theorem 13 (M. Pentus, 1996). A sequent $\rightarrow \Gamma$ is derivable in MCLL if and only if there exists a proof net for $\Gamma$.

Theorem 14. A sequent $\rightarrow \Gamma$ is derivable in MCLL' if and only if there exists a strong proof net for $\Gamma$.

Theorem 13 is proved in [24]. Here we modify this proof in order to prove Theorem 14.

Example 1.3. The following figure shows a proof net for the sequent Example 1.2:


Here (and further in the figures illustrating fragments of strong proof nets) the graphs $\mathcal{A}$ and $\mathcal{E}$ are drawn in the upper semiplane, and the relation $\prec_{\Gamma}$ (restricted to $\left.\Omega_{\Gamma}^{\otimes 8 \diamond}\right)$ is drawn in the lower semiplane.

In this example $\Omega_{\Gamma}=\left\{\langle\diamond, 1\rangle,\left\langle\bar{p}_{4}, 2\right\rangle,\langle\otimes, 3\rangle,\left\langle\bar{p}_{3}, 4\right\rangle,\langle 8,5\rangle,\left\langle p_{2}, 6\right\rangle,\langle\diamond, 7\rangle,\left\langle\bar{p}_{2}, 8\right\rangle\right.$, $\left.\langle\otimes, 9\rangle,\left\langle p_{1}, 10\right\rangle,\langle\diamond, 11\rangle,\left\langle\bar{p}_{1}, 12\right\rangle,\langle\otimes, 13\rangle,\left\langle p_{3}, 14\right\rangle,\langle 8,15\rangle,\left\langle p_{4}, 16\right\rangle\right\}$.

Example 1.4. The sequent $\rightarrow\left(\bar{p}_{2} \otimes\left(\bar{p}_{1} \ngtr p_{1}\right)\right) p_{2}$ (which is the translation of $\left(p_{1} \backslash p_{1}\right) \backslash p_{2} \rightarrow$ $p_{2}$, and the latter is derivable in $\mathrm{L}^{*}$, but not in L ) is derivable in MCLL, but not in MCLL'. The proof net below is not strong:


Indeed, let $X$ be the occurrence of the subformula $\bar{p}_{1} \otimes p_{1}$. Then $1(X)=\langle\otimes, 3\rangle$, $\mathrm{r}(X)=\langle\diamond, 7\rangle$ and $\widetilde{\mathcal{A}}(\mathrm{l}(X))=\widetilde{\mathcal{A}}(\mathrm{r}(X))=\langle\diamond, 7\rangle$, and this violates Axiom 6.

Example 1.5. There also exist sequents that have both a strong proof net and a proof net that is not strong. One example is the sequent $(p /(p \backslash p))(p \backslash p) \rightarrow p$. Its standard translation is $\rightarrow(\bar{p} \otimes p)((\bar{p} \ngtr p) \otimes \bar{p}) p$, and this sequent has two proof nets:

and


The first one is strong, and the second one is not. These two proof nets correspond to two Lambek derivations of the original sequent: $(p /(p \backslash p))(p \backslash p) \rightarrow p$ :

$$
\begin{aligned}
& \frac{p \rightarrow p \quad p \rightarrow p}{\frac{p(p \backslash p) \rightarrow p}{p \backslash p \rightarrow p \backslash p}} \begin{array}{l}
p \rightarrow p \\
(p /(p \backslash p))(p \backslash p) \rightarrow p
\end{array} \quad \text { and } \quad \frac{\frac{p \rightarrow p}{\rightarrow p \backslash p} \quad \frac{p \rightarrow p \quad p \rightarrow p}{p(p \backslash p) \rightarrow p}}{(p /(p \backslash p))(p \backslash p) \rightarrow p}
\end{aligned}
$$

The second derivation is valid only in $L^{*}$.
Lemma 1.6. If there exists a strong proof net $\mathfrak{N}=\left\langle\boldsymbol{\Omega}_{\Gamma}, \mathcal{A}, \mathcal{E}\right\rangle$, then $\Gamma$ contains at least two formulae.

Proof. Suppose the contrary: $\Gamma=A_{1}$. Let $X$ be the occurrence of $A_{1}$ (as a subformula of $\Gamma$ ). We have $1(X)=\mathrm{r}(X)=\langle\diamond, 0\rangle$, whence $\widetilde{\mathcal{A}}(\mathrm{l}(X))=\mathrm{l}(X)=\mathrm{r}(X)=\widetilde{\mathcal{A}}(\mathrm{r}(X))$. Contradiction with axiom 6.

Proof of Theorem 14. The "only if" part is proved by constructing the strong proof net inductively from the derivation of $\rightarrow \Gamma$ (exactly as in the proof of Theorem 13 [24]).

To prove the "if" part we proceed by induction on the number of connective occurrences in $\Omega_{\Gamma}$.

The induction base is trivial: in the case $\Omega_{\Gamma}^{8} \cup \Omega_{\Gamma}^{\otimes}=\emptyset$ our sequent can be either $\rightarrow p_{i} \bar{p}_{i}$ (axiom) or $\rightarrow \bar{p}_{i} p_{i}$ (derivable from $\rightarrow p_{i} \bar{p}_{i}$ by one application of (rot)).

Induction step. We define a new binary relation $\ll$ as the restriction of $\prec_{\Gamma} \cup \mathcal{A}$ to $\Omega_{\Gamma}^{\otimes} \cup \Omega_{\Gamma}^{\otimes}$. Due to the acyclicity of $\prec_{\Gamma} \cup \mathcal{A}$ there exists an element $\gamma \in \Omega_{\Gamma}^{8} \cup \Omega_{\Gamma}^{\otimes}$ that is maximal with respect to $\ll$. If $\gamma \in \Omega_{\Gamma}^{8}$, then we replace it by $\diamond$ (thus getting a strong proof net with fewer connective occurrences), use the induction hypothesis, and apply the rule $(\rightarrow>)$. The restriction of this rule is satisfied due to Lemma 1.6.

Now let $\gamma \in \Omega_{\Gamma}^{\otimes}, \beta=\mathcal{A}(\gamma)$. We have $\beta \in \Omega_{\Gamma}^{\diamond}$, because $\gamma$ is maximal with respect to $\ll$. We can assume that $\beta=\langle\diamond, 1\rangle$ (that is, $\beta$ is the leftmost occurrence of $\diamond$ ): in the other case we apply (rot) and do a cyclic transformation of the net.

We have $\Gamma=\Phi(A \otimes B) \Psi$, and the arc $\langle\gamma, \beta\rangle \in \mathcal{A}$ leads from the occurrence of $\otimes$ to the leftmost occurrence of $\diamond$. This arc divides the upper semiplane into two parts. We take the upper part as the strong proof net for $\rightarrow B \Psi$ and the lower part for $\rightarrow \Phi A$. Axioms 2, 3, 4, 5, and 6 are checked trivially. Axiom 1 is proved in [24], Lemma 7.10.

Now, using the induction hypothesis for $\rightarrow \Phi A$ and $\rightarrow B \Psi$, we conclude that these sequents are derivable in MCLL'. Therefore MCLL' $\vdash \rightarrow \Phi(A \otimes B) \Psi$ by application of $(\rightarrow \otimes)$.

We shall consider $\mathcal{E}$ a non-oriented graph on $\Omega_{\Gamma}$; the edges of $\mathcal{E}$ will be called links. Intuitively, links connect occurrences of atoms that come from one axiom leaf in the derivation tree. For a link $\mathbf{C}$ with vertices $\alpha$ and $\beta$ we define $\operatorname{In}(\mathbf{C}) \leftrightharpoons \operatorname{In}(\alpha, \beta)$ and $\operatorname{Out}(\mathbf{C}) \leftrightharpoons \operatorname{Out}(\alpha, \beta)$. The graph $\left\langle\Omega_{\Gamma}, \mathcal{E}\right\rangle$ is $<_{\Gamma}$-planar, whence $\#(\operatorname{In}(\mathbf{C}))=\#(\operatorname{Out}(\mathbf{C}))=$ 0 . Links divide the upper semiplane into regions. For each link $\mathbf{C}$ we define its inner and outer regions as the regions that have $\mathbf{C}$ as a part of their borders and intersect with $\operatorname{In}(\mathbf{C})$ and $\operatorname{Out}(\mathbf{C})$ respectively. Evidently, each region must contain at least one occurrence of $\ngtr$ or $\diamond$ (because $\mathcal{A} \cup \mathcal{E}$ is $<_{\Gamma}$-planar). On the other hand, it is easy to see that $\left|\Omega_{\Gamma}^{8 \diamond}\right|$ is equal to the number of regions, so each region contains only one occurrence of 8 or $\diamond$.

Let $X$ and $Y$ be two subformula occurrences such that $X \cap Y=\emptyset$. We define the fragment from $X$ to $Y$ as $\left\{\alpha \in \Omega_{\Gamma} \mid X<_{\Gamma}\{\alpha\}<_{\Gamma} Y\right\}$ if $X<_{\Gamma} Y$ and as $\left\{\alpha \in \Omega_{\Gamma} \mid\right.$ $\{\alpha\}<_{\Gamma} Y$ or $\left.X<_{\Gamma}\{\alpha\}\right\}$ if $Y<_{\Gamma} X$ (other cases are impossible).

If $\mathbf{C}$ is a link and $\mathscr{K}$ is a subset of $\Omega_{\Gamma}$, we define $\mathscr{D}(\mathbf{C}, \mathscr{K})=\operatorname{In}(\mathbf{C})$, if $\mathscr{K} \subseteq \operatorname{In}(\mathbf{C})$ and $\mathscr{D}(\mathbf{C}, \mathscr{K})=\operatorname{Out}(\mathbf{C})$ otherwise (if $\mathscr{K}$ is a fragment from one subformula occurrence to another, and it doesn't contain vertices of $\mathbf{C}$, then in this case we have $\mathscr{K} \subseteq \operatorname{Out}(\mathbf{C})$, thus getting $\mathscr{K} \subseteq \mathscr{D}(\mathbf{C}, \mathscr{K})$ always $)$. It is easy to see that $\#(\mathscr{D}(\mathbf{C}, \mathscr{K}))=0$.

Further we shall sometimes omit the word "occurrence".

## $1.8 L^{*}(\backslash)$-grammars for Context-Free Languages with the Empty Word

Theorem 15. Every context-free language is generated by some $L^{*}(\backslash)$-grammar.
To finish the proof of this theorem (see Section 1.4) is is sufficient to handle the case when the language contains the empty word. Let $M$ be context-free and $\epsilon \in M$. Using Theorem 5, we build an $\mathrm{L}(\backslash)$-grammar $\mathcal{G}$ for the language $M-\{\epsilon\}$.

Consider the type $D \leftrightharpoons((r \backslash r) \backslash((s \backslash s) \backslash q)) \backslash q$, where $q, r$, and $s$ are primitive types not occurring in grammar $\mathcal{G}$.

## Lemma 1.7.

1. $\mathrm{L}^{*}(\backslash) \vdash \rightarrow D$.
2. If a sequent $\Gamma \rightarrow p_{1}$, where $\Pi \neq \Lambda$, is a $G$-sequent and does not contain primitive types $q, r$, $s$, then the following equivalence holds:

$$
\mathrm{L}(\backslash) \vdash \Pi \rightarrow p_{1} \Longleftrightarrow \mathrm{~L}^{*}(\backslash) \vdash \Pi\left[p_{1}:=D\right] \rightarrow D .
$$

Proof. The first statement is established by the following derivation:

$$
\frac{\frac{r \rightarrow r}{\rightarrow r \backslash r}}{\frac{\frac{s \rightarrow s}{\rightarrow s \backslash s} q \rightarrow q}{(s \backslash s) \backslash q \rightarrow q}} \underset{\rightarrow((r \backslash r) \backslash((s \backslash s) \backslash q) \rightarrow q}{((s \backslash s) \backslash q)) \backslash q}
$$

The left-to-right implication in the second statement follows from the fact that any $\mathrm{L}(\backslash)$-derivable sequent is $\mathrm{L}^{*}(\backslash)$-derivable and the substitution rule.

The non-trivial part is the right-to-left implication in the second statement. We have $\mathrm{L}^{*}(\backslash) \vdash \Pi\left[p_{1}:=D\right] \rightarrow D$ and we need to prove that $\mathrm{L}(\backslash) \vdash \Pi \rightarrow p_{1}$. Due to Proposition 1.4, since $\Pi \rightarrow p_{1}$ is a G-sequent, it is sufficient to prove that $L^{*}(\backslash) \vdash \Pi \rightarrow p_{1}$.

Consider the translation to MCLL. It is easy to see that

$$
\widehat{D}=(\bar{q} \otimes(\bar{s} \ngtr s) \otimes(\bar{r} \ngtr r)) \ngtr q \quad \text { and } \quad \widehat{D}^{\perp}=\bar{q} \otimes((\bar{r} \otimes r) \ngtr(\bar{s} \otimes s) \ngtr q) .
$$

Let $\Pi=B_{1} \ldots B_{n}$ and let $\tilde{\Pi}={\widehat{B_{n}}}_{\tilde{\sim}}^{\perp} \ldots{\widehat{B_{1}}}^{\perp}$. Now we have MCLL $\vdash \rightarrow \tilde{\Pi}\left[p_{1}:=\widehat{D}\right] \widehat{D}$ and we need to prove that MCLL $\vdash \rightarrow \tilde{\Pi} p_{1}$.

In the sequent $\rightarrow \tilde{\Pi}\left[p_{1}:=\widehat{D}\right] \widehat{D}$ variables $q, r$, and $s$ appear only in the substituted $\widehat{D}$ and $\widehat{D}^{\perp}$ subformulae.

Consider the proof net for $\rightarrow \tilde{\Pi}\left[p_{1}:=\widehat{D}\right] \widehat{D}$. The occurrence of $\bar{s}$ from $\widehat{D}^{\perp}$ cannot be connected to an occurrence of $s$ from some $\widehat{D}^{\perp}$ (neither the same, nor some other), since in this case there would be two $\gamma \mathrm{s}$ in one area:

$$
\begin{aligned}
& \ldots \quad \bar{q} \otimes \bar{r} \otimes r \ngtr \bar{s} \otimes s \ngtr q \quad \ldots \quad \bar{q} \otimes \bar{r} \otimes r \ngtr \bar{s} \otimes s \ngtr q \quad \ldots \\
& \ldots \quad \bar{q} \otimes \bar{r} \otimes r \ngtr \bar{s} \otimes s \ngtr q \quad \ldots \quad \bar{q} \otimes \bar{r} \otimes r \ngtr \bar{s} \otimes s \ngtr q \quad \ldots \\
& \ldots \quad \bar{q} \otimes \bar{r} \otimes r \gg \bar{s} \otimes s \ngtr q \quad \ldots
\end{aligned}
$$

Hence this occurrence of $\bar{s}$ is connected to an $s$ from a $\widehat{D}$ occurrence.
The sequent $\Pi \rightarrow p_{1}$ is a G-sequent, whence there is only one positive occurrence of $p_{1}$. Therefore $\rightarrow \tilde{\Pi}\left[p_{1}:=\widehat{D}\right] \widehat{D}$ contains only one occurrence of $\widehat{D}$ (namely, the last formula) and, due to the previous argument, at most one occurrence of $\widehat{D}^{\perp}$. Consider two cases.

Case 1: there are no occurrences of $\widehat{D}^{\perp}$. Then $q$ and $\bar{q}$ from $\widehat{D}$ are connected, and $\tilde{\Pi}=\Lambda$ (in the other case there would be two $\diamond$ s in one area). Contradiction: $\Pi \neq \Lambda$.

Case 2: there is one occurrence of $\widehat{D}^{\perp}$. Then occurrences of $q, r, s, \bar{q}, \bar{r}$, and $\bar{s}$ from $\widehat{D}$ and $\widehat{D}^{\perp}$ are connected pairwise, and we can replace $\widehat{D}$ and $\widehat{D}^{\perp}$ by $p_{1}$ and $\bar{p}_{1}$, thus getting a proof net for $\rightarrow \tilde{\Pi} p_{1}$.

From this lemma it easily follows that substituting $D$ for $p_{1}$ in $\mathcal{G}$ yields an $L^{*}(\backslash)$ grammar generating the language $\mathfrak{L}(\mathcal{G}) \cup\{\epsilon\}=M$.

## Chapter 2

## The Lambek Calculus with the Unit

### 2.1 An Alternative Axiomatisation of $L_{1}$

We shall use an alternative axiomatisation of $\mathrm{L}_{\mathbf{1}}$, equivalent to the original one. The rule $(\mathbf{1} \rightarrow)$ can be considered a special case of the weakening rule. Every $\mathrm{L}_{1}$-derivation can be rebuilt in such a way that all applications of this rule will immediately follow the axioms. Formally this is done in the following way. Consider the calculus obtained from $\mathrm{L}_{1}$ by removing $(\mathbf{1} \rightarrow)$ and adding two new series of axioms: $\mathbf{1}^{k} \rightarrow \mathbf{1}(k \geqslant 0)$ and $\mathbf{1}^{k} p_{i} \mathbf{1}^{m} \rightarrow p_{i}$ $(k, m \geqslant 0, i \geqslant 1)$; denote them by $(\rightarrow \mathbf{1})_{\mathrm{w}}$ and $(\mathrm{ax})_{\mathrm{w}}$ respectively. The new calculus will be temporarily called $\widetilde{\mathrm{L}}_{1}$.

Proposition 2.1. For any sequent $\Pi \rightarrow B$, where $\Pi=A_{1} \ldots A_{n}, A_{1}, \ldots, A_{n}, B \in \mathrm{Tp}_{1}$, the following holds:

$$
\mathrm{L}_{\mathbf{1}} \vdash \Pi \rightarrow B \Longleftrightarrow \widetilde{\mathrm{~L}}_{1} \vdash \Pi \rightarrow B .
$$

Proof. To prove the right-to-left implication it is sufficient to establish derivability of the new axioms $(\rightarrow \mathbf{1})_{\mathrm{w}}(\mathrm{ax})_{\mathrm{w}}$ in $\mathrm{L}_{\mathbf{1}}$. This is done by induction on $k$ and $k+m$ respectively: the base case $(k=0$ and $k+m=0)$ corresponds to the axioms $(\rightarrow \mathbf{1})$ and $(\mathrm{ax})$ of $\mathrm{L}_{\mathbf{1}}$, and the induction step is an application of $(\mathbf{1} \boldsymbol{\rightarrow})$.

The left-to-right implication is established by induction on the length of the $\mathrm{L}_{1}$ derivation: if the rule $(\mathbf{1} \rightarrow$ ) was applied (possibly, several times) after an application of another rule, then the derivation can be rebuilt in such a way that the applications of $(\mathbf{1} \rightarrow)$ go, vice versa, before the application of the other rule. Applying $(\mathbf{1} \rightarrow)$ (several times) to the axioms (ax) and $(\rightarrow \mathbf{1})$ yields $(\mathrm{ax})_{\mathrm{w}}$ and $(\rightarrow \mathbf{1})_{\mathrm{w}}$ respectively.

Further we shall use this new axiomatisation for $\mathrm{L}_{\mathbf{1}}$.

### 2.2 Reduction of $L_{1}$ to $L^{*} . L_{1}$-grammars

We present a substitution that reduces derivability in $\mathrm{L}_{\mathbf{1}}$ to derivability in $\mathrm{L}^{*}$.
Theorem 16. For any sequent $\Pi \rightarrow C$ built from types that belong to $\mathrm{Tp}_{\mathbf{1}}$ and for any primitive type $q$ not occurring in $\Pi \rightarrow C$, the following equivalence holds:

$$
\mathrm{L}_{\mathbf{1}} \vdash \Pi \rightarrow C \Longleftrightarrow \mathrm{~L}^{*} \vdash(\Pi \rightarrow C)\left[p_{i}:=\left(\mathbf{1} \cdot p_{i}\right) \cdot \mathbf{1}\right][\mathbf{1}:=q \backslash q] .
$$

Here and further the shorthand " $p_{i}:=\left(\mathbf{1} \cdot p_{i}\right) \cdot p_{i}$ " means that the substitution is performed for every $i$.

Lemma 2.2. For every $k \geqslant 0$ the following holds: $\mathrm{L}^{*}(\backslash) \vdash(q \backslash q)^{k} \rightarrow q \backslash q$.
Proof.

$$
\begin{array}{r}
\frac{q \rightarrow q}{\frac{q \rightarrow q \quad q \rightarrow q}{q(q \backslash q) \rightarrow q}} \\
\frac{q}{q \rightarrow q} \frac{q \rightarrow q \quad \frac{q(q \backslash q) \ldots(q \backslash q) \rightarrow q}{q(q \backslash q)(q \backslash q) \ldots(q \backslash q) \rightarrow q}}{(q \backslash q)(q \backslash q)(q \backslash q) \ldots(q \backslash q) \rightarrow q} \\
(q \backslash q)(q \backslash q)(q \backslash q) \ldots(q \backslash q) \rightarrow q \backslash q
\end{array}
$$

Consider an auxiliary calculus $\mathrm{L}_{\mathbf{1}}^{-}$, which is obtained from $\mathrm{L}^{*}$ by adding axioms $(\rightarrow \mathbf{1})_{\mathrm{w}}$. It is clear that every $\mathrm{L}_{\mathbf{1}}^{-}$-derivable sequent is $\mathrm{L}_{\mathbf{1}}$-derivable.

Lemma 2.3. For every sequent $\Pi \rightarrow C$ built from types that belong to $\mathrm{Tp}_{\mathbf{1}}$ the following equivalences hold:

$$
\begin{aligned}
\mathrm{L}_{\mathbf{1}} \vdash \Pi \rightarrow C \Longleftrightarrow \mathrm{~L}_{\mathbf{1}} \vdash(\Pi \rightarrow C)\left[p_{i}:=\left(\mathbf{1} \cdot p_{i}\right) \cdot \mathbf{1}\right] & \Longleftrightarrow \\
& \Longleftrightarrow \mathrm{L}_{\mathbf{1}}^{-} \vdash(\Pi \rightarrow C)\left[p_{i}:=\left(\mathbf{1} \cdot p_{i}\right) \cdot \mathbf{1}\right] .
\end{aligned}
$$

Proof. The first equivalence follows from the fact that $p_{i} \leftrightarrow_{\mathrm{L}_{1}}\left(\mathbf{1} \cdot p_{i}\right) \cdot \mathbf{1}$.
In the second equivalence the right-to-left implication is obvious. Let us prove the other one: we shall deduce the third statement from the first one (which is equivalent to the second one). We substitute $\left(\mathbf{1} \cdot p_{i}\right) \cdot \mathbf{1}$ in the $\mathrm{L}_{\mathbf{1}}$-derivation of $\Pi \rightarrow C$. It is easy to see that this substitution conserves the $(\rightarrow \mathbf{1})_{\mathrm{w}}$ axioms and all rules. Now it is sufficient to check that the result of such a substitution in $(\mathrm{ax})_{\mathrm{w}}$ is derivable in $\mathrm{L}_{\mathbf{1}}^{-}$:

$$
\begin{aligned}
& \frac{\mathbf{1}^{k+1} \rightarrow \mathbf{1} \quad p_{i} \rightarrow p_{i}}{\mathbf{1}^{k+1} p_{i} \rightarrow \mathbf{1} \cdot p_{i}} \quad \mathbf{1}^{m+1} \rightarrow \mathbf{1} \\
& \frac{\mathbf{1}^{k} \mathbf{1} p_{i} \mathbf{1}^{m+1} \rightarrow\left(\mathbf{1} \cdot p_{i}\right) \cdot \mathbf{1}}{\mathbf{1}^{k}\left(\mathbf{1} \cdot p_{i}\right) \mathbf{1} \mathbf{1}^{m} \rightarrow\left(\mathbf{1} \cdot p_{i}\right) \cdot \mathbf{1}} \\
& \mathbf{1}^{k}\left(\left(\mathbf{1} \cdot p_{i}\right) \cdot \mathbf{1}\right) \mathbf{1}^{m} \rightarrow\left(\mathbf{1} \cdot p_{i}\right) \cdot \mathbf{1}
\end{aligned}
$$

The next two lemmas essentially repeat the argument from [19] about the closed (without variables but with constants) fragment of multiplicative cyclic linear logic.

Lemma 2.4. If $\mathrm{L}_{\mathbf{1}}^{-} \vdash \Pi \rightarrow C \quad q \in \operatorname{Pr}$, then $\mathrm{L}^{*} \vdash(\Pi \rightarrow C)[\mathbf{1}:=q \backslash q]$.
Proof. Perform the substitution in the $\mathrm{L}_{1}^{-}$-derivation of $\Pi \rightarrow C$. Axioms $(\mathrm{ax})_{\mathrm{w}}$ and rules of inference will remain untouched. Axioms $(\rightarrow \mathbf{1})_{\mathrm{w}}$ will transform into sequents $(q \backslash q)^{k} \rightarrow q \backslash q$, which are derivable in L* by Lemma 2.2.

Lemma 2.5. If $\mathrm{L}^{*} \vdash(\Pi \rightarrow C)[\mathbf{1}:=q \backslash q]$ and $q$ is a primitive type that does not occur in $\Pi \rightarrow C$, then $\mathrm{L}_{\mathbf{1}} \vdash \Pi \rightarrow C$.

Proof. Let L* ${ }^{*} \vdash(\Pi \rightarrow C)[\mathbf{1}:=q \backslash q]$. Consider the sequent $(\Pi \rightarrow C)[\mathbf{1}:=q \backslash q][q:=\mathbf{1}]$. On the one hand, it is derivable in $\mathrm{L}_{\mathbf{1}}$, since $(\Pi \rightarrow C)[\mathbf{1}:=q \backslash q]$ is derivable in $\mathrm{L}_{\mathbf{1}}$ (due to the conservativity of $\mathrm{L}_{\mathbf{1}}$ over $\mathrm{L}^{*}$ ) and the substitution rule is valid in $\mathrm{L}_{\mathbf{1}}$. On the other hand, the sequent involved is actually $(\Pi \rightarrow C)[\mathbf{1}:=\mathbf{1} \backslash \mathbf{1}]$, because occurrences of $q$ could appear only inside the types $q \backslash q$ that are substituted for $\mathbf{1}$. Therefore the derivability of this sequent in $\mathrm{L}_{\mathbf{1}}$ is equivalent to the derivability of $\Pi \rightarrow C$ (since $\mathbf{1} \leftrightarrow \mathbf{1} \backslash \mathbf{1}$ ).

Proof of Theorem 16.

$$
\begin{aligned}
\mathrm{L}_{\mathbf{1}} \vdash \Pi \rightarrow C \Longrightarrow & \mathrm{~L}_{\mathbf{1}}^{-} \vdash(\Pi \rightarrow C)\left[p_{i}:=\left(\mathbf{1} \cdot p_{i}\right) \cdot \mathbf{1}\right] \Longrightarrow \\
& \Longrightarrow \mathrm{L}^{*} \vdash(\Pi \rightarrow C)\left[p_{i}:=\left(\mathbf{1} \cdot p_{i}\right) \cdot \mathbf{1}\right][\mathbf{1}:=q \backslash q] \Longrightarrow \\
& \Longrightarrow \mathrm{L}_{\mathbf{1}} \vdash(\Pi \rightarrow C)\left[p_{1}:=\left(\mathbf{1} \cdot p_{i}\right) \cdot \mathbf{1}\right] \Longrightarrow \mathrm{L}_{\mathbf{1}} \vdash \Pi \rightarrow C .
\end{aligned}
$$

Here the first and the fourth implications hold due to Lemma 2.3, the second one holds due to Lemma 2.4, and the third one holds due to Lemma 2.5.

Theorem 16 easily yields a corollary concerning $\mathrm{L}_{1}$-grammars. (In [30] it is stated that the theorem below can be proved by the same means as for $L^{*}$, but this is not true.)

Theorem 17. The class of languages generated by $\mathrm{L}_{\mathbf{1}}$-grammars coincides with the class of all context-free languages.

Proof. Let $\mathcal{G}$ be an $\mathrm{L}_{\mathbf{1}}$-grammar. Then, due to Theorem 16, the $\mathrm{L}^{*}$-grammar $\mathcal{G}\left[p_{i}:=\right.$ $\left.\left(\mathbf{1} \cdot p_{i}\right) \cdot \mathbf{1}\right][\mathbf{1}:=q \backslash q]$ generates the same language as $\mathcal{G}$. On the other hand, this language is context-free by Theorem 4 .

The inverse inclusion is due to conservativity of $\mathrm{L}_{\mathbf{1}}$ over $\mathrm{L}^{*}$ and Theorem 15.

### 2.3 Reduction of $\mathrm{L}_{1}(\backslash)$ to $\mathrm{L}^{*}(\backslash)$. A Polynomial Algorithm for the Derivability Problem in $\mathrm{L}_{1}(\backslash)$

Now we shall need a refinement of Theorem 16, a substitution with only one division.
Theorem 18. For any sequent $\Pi \rightarrow C$ built from types belonging to $\mathrm{Tp}_{\mathbf{1}}(\backslash)$ and for any primitive type $q$ not occurring in $\Pi \rightarrow C$, the following equivalence holds:

$$
\mathrm{L}_{\mathbf{1}}(\backslash) \vdash \Pi \rightarrow C \Longleftrightarrow \mathrm{~L}^{*}(\backslash) \vdash(\Pi \rightarrow C)\left[\mathbf{1}:=q \backslash q, p_{i}:=(q \backslash q) \backslash\left(p_{i} \backslash q\right)\right] .
$$

This theorem yields polynomial decidability of the derivability problem in $\mathrm{L}_{\mathbf{1}}(\backslash)$ (derivability problems for $\mathrm{L}_{\mathbf{1}}, \mathrm{L}_{\mathbf{1}}(\backslash, /)$, and $\mathrm{L}_{\mathbf{1}}(\cdot, \backslash)$ are NP-complete due to conservativity and Theorem 8):

Theorem 19. There exists an algorithm that checks derivability of a sequent in $\mathrm{L}_{\mathbf{1}}(\backslash)$ in $O\left(n^{3}\right)$ steps, where $n$ is the length of the sequent.

Proof. The algorithm is as follows: first we build the sequent $(\Pi \rightarrow C)\left[\mathbf{1}:=q \backslash q, p_{i}:=\right.$ $(q \backslash q) \backslash\left(p_{i} \backslash q\right)$ ], where $q$ is a new primitive type. The length of the sequent grows not more than seven times. Then we use Savateev's algorithm (Theorem 9) to check derivability of this sequent in $L^{*}(\backslash)$.

Before proving Theorem 18, we prove several lemmas.

Lemma 2.6. For any $k, m \geqslant 0, p_{i} \in \operatorname{Pr}$ the following holds:

$$
\mathrm{L}^{*}(\backslash) \vdash(q \backslash q)^{k}\left((q \backslash q) \backslash\left(p_{i} \backslash q\right)\right)(q \backslash q)^{m} \rightarrow(q \backslash q) \backslash\left(p_{i} \backslash q\right) .
$$

Proof.

$$
\begin{array}{r}
\frac{(q \backslash q)^{k+1} \rightarrow q \backslash q}{\frac{p_{i} \rightarrow p_{i} \quad q(q \backslash q)^{m} \rightarrow q}{p_{i}\left(p_{i} \backslash q\right)(q \backslash q)^{m} \rightarrow q}}\left(\underline{\left(p_{i} \backslash q\right)(q \backslash q)^{m} \rightarrow p_{i} \backslash q}\right. \\
(q \backslash q)^{k+1}\left((q \backslash q) \backslash\left(p_{i} \backslash q\right)\right)(q \backslash q)^{m} \rightarrow(q \backslash q) \backslash\left(p_{i} \backslash q\right)
\end{array}
$$

Sequents $(q \backslash q)^{k+1} \rightarrow q \backslash q$ and $q(q \backslash q)^{m} \rightarrow q$ were derived in the proof of Lemma 2.2.
Lemma 2.7. If $\mathrm{MCLL}_{\mathbf{1}, \perp} \vdash \rightarrow \Gamma\left[p_{i}:=\bar{p}_{i}\right]$, then $\mathrm{MCLL}_{\mathbf{1}, \perp} \vdash \rightarrow \Gamma$.
Proof. After substituting $\bar{p}_{i}$ for $p_{i}$ in the $\mathrm{MCLL}_{\mathbf{1}, \perp}$-derivation of a sequent, the rules of inference and axiom $\rightarrow \mathbf{1}$ remain valid, and axioms $\rightarrow p_{i} \bar{p}_{i}$ become sequents $\rightarrow \bar{p}_{i} p_{i}$, derivable by one application of (rot). Finally, notice that $\Gamma\left[p_{i}:=\bar{p}_{i}\right]\left[p_{i}:=\bar{p}_{i}\right]=\Gamma$.

Lemma 2.8. The following equivalences hold in $\mathrm{MCLL}_{\mathbf{1}, \perp}: \mathbf{1} \ngtr \perp \leftrightarrow \mathbf{1}$ and $(\mathbf{1} \otimes \perp) \ngtr\left(\bar{p}_{i} \ngtr \perp\right) \leftrightarrow$ $\bar{p}_{i}$.

Proof. Since $A \leftrightarrow_{\mathrm{MCLL}_{1, \perp}} B$ means that $\mathrm{MCLL}_{\mathbf{1}, \perp} \vdash \rightarrow A^{\perp} B \quad \mathrm{MCLL}_{\mathbf{1}, \perp} \vdash \rightarrow B^{\perp} A$, and $(\mathbf{1} \diamond \perp)^{\perp}=\mathbf{1} \otimes \perp, \mathbf{1}^{\perp}=\perp, \bar{p}_{i}^{\perp}=p_{i},\left((\mathbf{1} \otimes \perp) \ngtr\left(\bar{p}_{i} \oslash \perp\right)\right)^{\perp}=\left(\mathbf{1} \otimes p_{i}\right) \otimes(\mathbf{1} \oslash \perp)$, we need to derive the following sequents in $\mathrm{MCLL}_{1, \perp}: \rightarrow(\mathbf{1} \otimes \perp) \mathbf{1}, \rightarrow \perp(\mathbf{1} \ngtr \perp), \rightarrow$ $\left(\left(\mathbf{1} \otimes p_{i}\right) \otimes(\mathbf{1} \oslash \perp)\right) \bar{p}_{i} \rightarrow p_{i}\left((\mathbf{1} \otimes \perp) \ngtr\left(\bar{p}_{i} \ngtr \perp\right)\right)$. The derivations are as follows:

Proof of Theorem 18. Prove the left-to-right implication. Consider an $\mathrm{L}_{\mathbf{1}}(\backslash)$-derivation of $\Pi \rightarrow C$. Perform the substitution $1:=q \backslash q, p_{i}:=(q \backslash q) \backslash\left(p_{i} \backslash q\right)$ in this derivation. The rules of inference remain valid, axioms $(\rightarrow \mathbf{1})_{\mathrm{w}}$ transform to sequents $(q \backslash q)^{k} \rightarrow$ $q \backslash q$, derivable by Lemma 2.2, and axioms (ax) ${ }_{\mathrm{w}}$ transform to sequents $(q \backslash q)^{k} \rightarrow q \backslash q$, derivable by Lemma 2.6. Therefore, the sequent $(\Pi \rightarrow C)\left[\mathbf{1}:=q \backslash q, p_{i}:=(q \backslash q) \backslash\left(p_{i} \backslash q\right)\right]$ is derivable in $L^{*}(\backslash)$.

Now prove the right-to-left implication. Let $L^{*}(\backslash) \vdash(\Pi \rightarrow C)\left[1:=q \backslash q, p_{i}:=\right.$ $\left.(q \backslash q) \backslash\left(p_{i} \backslash q\right)\right]$. Then $\operatorname{MCLL}_{\mathbf{1}, \perp} \vdash(\widehat{\Pi} \rightarrow \widehat{C})\left[\mathbf{1}:=\bar{q} \ngtr q, p_{i}:=(\bar{q} \otimes q) \ngtr\left(\bar{p}_{i} \ngtr q\right)\right]$. Substitute $\perp$ for $q$. This yields to a sequent derivable in $\operatorname{MCLL}_{1, \perp}$, and, since $q$ does not occur
in $\widehat{\Pi} \rightarrow \widehat{C}$, this sequent is actually $(\widehat{\Pi} \rightarrow \widehat{C})\left[\mathbf{1}:=\mathbf{1} \ngtr \perp, p_{i}:=(\mathbf{1} \otimes \perp) \ngtr\left(\bar{p}_{i} \ngtr \perp\right)\right]$. Due to Lemma 2.8, $\mathrm{MCLL}_{\mathbf{1}, \perp} \vdash(\widehat{\Pi} \rightarrow \widehat{C})\left[p_{i}:=\bar{p}_{i}\right]$. Finally, by Lemma $2.7 \mathrm{MCLL}_{\mathbf{1}, \perp} \vdash \widehat{\Pi} \rightarrow \widehat{C}$, whence $\mathrm{L}_{\mathbf{1}}(\backslash) \vdash \Pi \rightarrow C$. Q. E. D.

## Chapter 3

## The Lambek Calculus with One Primitive Type

### 3.1 Reduction of $\mathrm{L}\left(\backslash ; p_{1}, \ldots, p_{N}\right)$ to $\mathrm{L}(\backslash ; p)$

Let $N$ be an arbitrary natural number. We build a faithful substitution that reduces derivability in $\mathrm{L}\left(\backslash ; p_{1}, \ldots, p_{N}\right)$ to derivability in $\mathrm{L}(\backslash ; p)$.

Let $p^{m} \leftrightharpoons \underbrace{p \cdot \ldots \cdot p}_{m \text { times }}$.
Now construct the types $A_{1}, \ldots, A_{N}$ :

$$
A_{k} \leftrightharpoons\left(p^{k+1} \cdot(((p \cdot p) \backslash p) \backslash p) \cdot p^{N-k+1}\right) \backslash p, \quad k=1, \ldots, N .
$$

Note that in $\mathrm{L}(\backslash)$ there is no multiplication; the • connective here (multiplication in the denominator) is just a shortcut used to minimise the number of brackets: $\left(D_{1} \ldots\right.$. $\left.D_{m}\right) \backslash C \leftrightharpoons D_{m} \backslash\left(D_{m-1} \backslash\left(D_{m-2} \backslash \ldots \backslash\left(D_{1} \backslash C\right) \ldots\right)\right)$. This notation is consistent with the "real" multiplication: $\left(D_{1} \cdot \ldots \cdot D_{m}\right) \backslash C \leftrightarrow_{\mathrm{L}(\cdot,))} D_{m} \backslash\left(D_{m-1} \backslash\left(D_{m-2} \backslash \ldots \backslash\left(D_{1} \backslash C\right) \ldots\right)\right.$ ).

Theorem 20. For any sequent $\Pi \rightarrow C$, where $\Pi=B_{1} \ldots B_{m}, m \geqslant 1$, and $B_{1}, \ldots, B_{m}, C \in \mathrm{Tp}\left(\backslash ; p_{1}, \ldots, p_{N}\right)$, the following holds:

$$
\mathrm{L}\left(\backslash ; p_{1}, \ldots, p_{N}\right) \vdash \Pi \rightarrow C \Longleftrightarrow \mathrm{~L}(\backslash ; p) \vdash(\Pi \rightarrow C)\left[p_{1}:=A_{1}, \ldots, p_{N}:=A_{N}\right] .
$$

Theorem 21. For every sequent $\Pi \rightarrow C$, where $\Pi=B_{1} \ldots B_{m}, m \geqslant 0$, and $B_{1}, \ldots, B_{m}, C \in \operatorname{Tp}\left(\backslash ; p_{1}, \ldots, p_{N}\right)$, the following holds:

$$
\mathrm{L}^{*}\left(\backslash ; p_{1}, \ldots, p_{N}\right) \vdash \Pi \rightarrow C \Longleftrightarrow \mathrm{~L}^{*}(\backslash ; p) \vdash(\Pi \rightarrow C)\left[p_{1}:=A_{1}, \ldots, p_{N}:=A_{N}\right] .
$$

Due to conservativity of $\mathrm{MCLL}^{\prime}$ over $\mathrm{L}(\backslash)$ (and, resp., $\operatorname{MCLL}^{\prime}(p)$ over $\mathrm{L}(\backslash ; p)$ and $\operatorname{MCLL}^{\prime}\left(p_{1}, \ldots, p_{N}\right)$ over L( $\left.\left.\backslash ; p_{1}, \ldots, p_{N}\right)\right)$ and MCLL over L* $(\backslash)($ resp., $\operatorname{MCLL}(p)$ over $\mathrm{L}^{*}(\backslash ; p)$ and $\operatorname{MCLL}\left(p_{1}, \ldots, p_{N}\right)$ over $\left.\mathrm{L}^{*}\left(\backslash ; p_{1}, \ldots, p_{N}\right)\right)$ these two theorems are corollaries from the following two stronger statements, which we are going to prove in the next section:

Theorem 22. For any sequent $\rightarrow \Gamma$, where $\Gamma=B_{1} \ldots B_{m}, m \geqslant 2$, and $B_{1}, \ldots, B_{m} \in$ $\operatorname{Fm}\left(p_{1}, \ldots, p_{N}\right)$, the following holds:

$$
\operatorname{MCLL}^{\prime}\left(p_{1}, \ldots, p_{N}\right) \vdash \rightarrow \Gamma \Longleftrightarrow \operatorname{MCLL}^{\prime}(p) \vdash \rightarrow \Gamma\left[p_{1}:=\widehat{A}_{1}, \ldots, p_{N}:=\widehat{A}_{N}\right] .
$$

Theorem 23. For every sequent $\rightarrow \Gamma$, where $\Gamma=B_{1} \ldots B_{m}, m \geqslant 1$, and $B_{1}, \ldots, B_{m} \in$ $\operatorname{Fm}\left(p_{1}, \ldots, p_{N}\right)$, the following holds:

$$
\operatorname{MCLL}\left(p_{1}, \ldots, p_{N}\right) \vdash \rightarrow \Gamma \Longleftrightarrow \operatorname{MCLL}(p) \vdash \rightarrow \Gamma\left[p_{1}:=\widehat{A}_{1}, \ldots, p_{N}:=\widehat{A}_{N}\right]
$$

A faithful substitution that reduces derivability in $\operatorname{MCLL}\left(p_{1}, \ldots, p_{N}\right)$ to derivability in $\operatorname{MCLL}(p)$ was earlier presented by Métayer [19]. The construction of Métayer can be easily modified to obtain a substitution that reduces derivability in $L^{*}\left(\backslash, / ; p_{1}, \ldots, p_{N}\right)$ to derivability in $\mathrm{L}^{*}(\backslash, / ; p)$ (we substitute $A_{k} \leftrightharpoons p^{k} \backslash p / p^{N+1-k}$ for $p_{k}$; here we use both divisions); a similar result can be obtained for L using strong proof nets.

Yet another substitution reducing derivability in $\mathrm{L}(\backslash)$ to derivability in $\mathrm{L}(\backslash ; p)$ was independently presented by Hendriks in [13]. Hendriks' construction uses only one division, but the types used in it have exponential size w. r. t. $N$. Therefore this construction cannot be used to prove NP-completeness of $\mathrm{L}(\cdot, \backslash ; p), \mathrm{L}(\backslash, / ; p)$ and their variants allowing empty antecedents. Also, in our opinion, Hendriks' proof is more complicated then the proof presented here.

### 3.2 Faithfullness of the Substitution (Proof)

We prove Theorems 20 and 21 in parallel. The left-to-right implications are trivial.
Let $R_{k} \leftrightharpoons \widehat{A}_{k}$. It is easy to see that

$$
\begin{aligned}
& R_{k}=\underbrace{\bar{p} \ngtr \ngtr \bar{p}}_{N-k+1} \ngtr\left(\bar{p} \otimes\left(\bar{p} \ngtr \bar{p} \ngtr p_{*}\right)\right) \gtrdot \underbrace{\bar{p} \ngtr \ldots \gtrdot \bar{p}}_{k+1} \ngtr p_{+} ; \\
& R_{k}^{\perp}=\bar{p}_{+} \otimes \underbrace{p \otimes \ldots \otimes p}_{k+1} \otimes\left(\left(\bar{p}_{*} \otimes p \otimes p\right) \gtrdot p\right) \otimes \underbrace{p \otimes \ldots \otimes p}_{N-k+1} .
\end{aligned}
$$

(We assume that 8 s associate to the right and $\otimes \mathrm{s}$ associate to the left.) The subscripts * and + here mark concrete occurrences of $p$ and $\bar{p}$ for further reference. Let us call $R_{k}$ a positive formula and $R_{k}^{\perp}$ a negative formula.

If $X$ is an occurrence of $R_{k}$, then $\#(X)=-(N+3)$, and if $Y$ is an occurrence of $R_{k}^{\perp}$, then $\#(Y)=N+3$. We shall use these count checks later.

The sequent $\rightarrow \Gamma\left[p_{1}:=R_{1}, \ldots, p_{N}:=R_{N}\right]$ is derivable in $\operatorname{MCLL}^{\prime}(p)$ (resp., in $\operatorname{MCLL}(p))$. Hence there exists a strong proof net (resp., just a proof net) $\mathfrak{N}$ for this sequent. We shall modify $\mathfrak{N}$ to get a proof net $\mathfrak{N}^{\prime}$ for $\rightarrow \Gamma$. If the original proof net were strong, $\mathfrak{N}^{\prime}$ will also be strong.

First we shall prove some lemmata about $\mathfrak{N}$.
Lemma 3.1. The number of positive formula occurrences is equal to the number of negative formula occurrences.

Proof. Suppose there are $m_{1}$ positive formula occurrences and $m_{2}$ negative formula occurrences. Then we have $0=\#\left(\Omega_{\Gamma}\right)=\left(m_{2}-m_{1}\right)(N+3)$, whence $m_{2}=m_{1}$.

Lemma 3.2. Any occurrence of $p_{*}$ from $R_{k}$ is connected to an occurrence of $\bar{p}_{*}$ from some $R_{k^{\prime}}^{\perp}$ (possibly, $k \neq k^{\prime}$ ).

Proof. Suppose the contrary. Let C be a link from some $p_{*}$ that does not lead to $\bar{p}_{*}$. We consider 3 cases:

Case 1: C leads to an occurrence of $\bar{p}$ from the same $R_{k}$. Since $\#(\operatorname{In}(\mathbf{C}))=0$, this is the neighbour occurrence of $\bar{p}$. But then there are two 8 in the outer region of $\mathbf{C}$. Contradiction.

Case 2: $\mathbf{C}$ leads to an occurrence of $\bar{p}$ from another $R_{k^{\prime}}$. There are $>$ connectives on both sides of $p_{*}$ and on at least one side of any $\bar{p}$ from $R_{k^{\prime}}$, therefore either in the inner or in the outer region of $\mathbf{C}$ there are two 8 connectives. Contradiction.

Case 3: C leads to $\bar{p}_{+}$(from some $R_{k^{\prime}}^{\perp}$ ):

$$
\bar{p} \ngtr \ldots \diamond \bar{p} \gtrdot \bar{p} \otimes \bar{p} \ngtr \bar{p} \ngtr p_{*} \ngtr \underbrace{\bar{p} \diamond \ldots \diamond \bar{p}}_{k+1} \gtrdot p_{+} \square \mathscr{K} \bar{p}_{+} \otimes p \otimes \ldots \otimes p \otimes \bar{p}_{*} \otimes p \otimes p \ngtr p \otimes p \otimes \ldots \otimes p
$$

Let $\mathscr{K}$ be the fragment from $R_{k}$ to $R_{k^{\prime}}^{\perp}$. The fragment $\mathscr{K}$ consists of several occurrences of positive and negative formulae and connectives between them. Let $m_{1}$ be the number of positive formulae there and $m_{2}$ be the number of the negative ones. Then we get $0=\#(\mathscr{D}(\mathbf{C}, \mathscr{K}))=1-(k+1)+\#(\mathscr{K})=-k+\left(m_{2}-m_{1}\right)(N+3)$, therefore $k=$ $\left(m_{2}-m_{1}\right)(N+3)$. This is absurd, because $1 \leq k \leq N$.

Lemma 3.3. Any occurrence of $\bar{p}_{*}$ from $R_{k}^{\perp}$ is connected to an occurrence of $p_{*}$ from some $R_{k^{\prime}}$.

Proof. Any occurrence of $p_{*}$ is connected to an occurrence of $\bar{p}_{*}$, and the numbers of $p_{*}$ and $\bar{p}_{*}$ occurrences coincide.

Lemma 3.4. If an occurrence of $p_{*}$ from $R_{k}$ is connected to an occurrence of $\bar{p}_{*}$ from $R_{k^{\prime}}^{\perp}$, then $k=k^{\prime}$.

Proof.

Let $\mathscr{K}$ be the fragment from $R_{k}$ to $R_{k^{\prime}}^{\perp}$. Occurrences of positive formulae inside $\mathscr{K}$ are in one-to-one correspondence with occurrences of negative formulae there (by the links connecting $p_{*}$ and $\left.\bar{p}_{*}\right)$, thus $\#(\mathscr{K})=0$. Therefore $0=\#(\mathscr{D}(\mathbf{C}, \mathscr{K}))=k^{\prime}-k+\#(\mathscr{K})$. Hence $k=k^{\prime}$.

Lemma 3.5. Any occurrence of $p_{+}$from $R_{k}$ is connected to an occurrence of $\bar{p}_{+}$from some $R_{k^{\prime}}^{\perp}$ (possibly, $k \neq k^{\prime}$ ).

Proof. We consider several cases:
Case 1: the occurrence of $p_{+}$is connected to an occurrence of $\bar{p}$ from the same $R_{k}$ by a link C. Since $\#(\operatorname{In}(\mathbf{C}))=0$, it is the rightmost occurrence of $\bar{p}$ :


But then immediately to the right of $R_{k}$ there is an occurrence $\tau$ of $\otimes$ (otherwise there would be two occurrences of $\varnothing$ in one region) connected by an $\mathcal{A}$-arc with the occurrence $\pi$ of $\ngtr$ on the left side of $\mathbf{C}$ (due to $<_{\Gamma}$-planarity of $\mathcal{A} \cup \mathcal{E}$ ). On the other hand, $\pi \prec_{\Gamma} \tau$ (since $\tau$ is situated immediately next to $R_{k}$, and therefore $R_{k}$ is a subformula
of the formula where $\tau$ is the occurrence of the main connective). Contradiction with the acyclicity of $\mathcal{A} \cup \prec_{\Gamma}$.

Case 2: $p_{+}$is connected to the occurrence of $\bar{p}$ in $R_{k^{\prime}}$, which is the third to the left from $p_{*}$, by a link $\mathbf{C}$ :

Let $\mathscr{K}$ be the fragment from $R_{k}$ to $R_{k^{\prime}}$. The same argument as in Lemma 3.4 shows that $\#(\mathscr{K})=0$. But then $\#(\mathscr{D}(\mathbf{C}, \mathscr{K}))=-\left(N-k^{\prime}+1\right)+\#(\mathscr{K}) \neq 0$. Contradiction.

Case 3: $p_{+}$is connected by a link $\mathbf{C}$ to an occurrence of $\bar{p}$ from another $R_{k^{\prime}}$, but not the third to the left from $p_{*}$. In this case either in the inner (if $R_{k^{\prime}}$ lies to the left from $R_{k}$ ) or in the outer region of $\mathbf{C}$ there are two occurrences of 8 . Contradiction.

Case 4: $p_{+}$is connected to $\bar{p}_{*}$ from some $R_{k^{\prime}}^{\perp}$. Contradiction with Lemma 3.3: $\bar{p}_{*}$ is connected to an occurrence of $p_{*}$, but not $p_{+}$.

So the only possible situation is the 5 th case: $p_{+}$is connected to an occurrence of $\bar{p}_{+}$from some $R_{k^{\prime}}^{\perp}$.

Lemma 3.6. Any occurrence of $\bar{p}_{+}$from $R_{k}^{\perp}$ is connected to an occurrence of $p_{+}$from some $R_{k^{\prime}}$.

Proof. The same argument as in Lemma 3.3.
Let us call an occurrence of a connective old if it is not inside an occurrence of a positive or negative formula (thus this occurrence comes from the original sequent $\rightarrow \Gamma$ ).

Lemma 3.7. If an occurrence $\tau$ of $\otimes$ is old, then $\mathcal{A}(\tau)$ is also old.
Proof. Suppose the contrary. Let $\mathcal{A}(\tau)$ be not old and let $\mathcal{A}(\tau)<_{\Gamma} \tau$ (in the other case we proceed symmetrically with respect to the $\operatorname{arc}\langle\tau, \mathcal{A}(\tau)\rangle)$. Consider several cases:

Case 1: $\mathcal{A}(\tau)$ lies inside some $R_{k}^{\perp}$ :


Let $\mathscr{D}=\operatorname{In}(\tau, \mathcal{A}(\tau))$ and $\mathscr{K}$ be the fragment of $\Omega_{\Gamma}$ between $R_{k}^{\perp}$ and $\tau$. We have $\#(\mathscr{K})=0$ and $\#(\mathscr{D})=0$. Contradiction.

Case 2: $\mathcal{A}(\tau)$ is an occurrence of $\varnothing$ inside $R_{k}$, but not the second from the right side. From the two sets $\operatorname{In}(\tau, \mathcal{A}(\tau))$ and $\operatorname{Out}(\tau, \mathcal{A}(\tau))$ we take the one not containing $p_{*}$ from this $R_{k}$ and call it $\mathscr{D}$. For $\mathscr{K}$ we take the subset of $\Omega_{\Gamma}$ containing all elements of $\mathscr{D}$, except those from $R_{k}$. Now we proceed exactly as in case 1 .

Case 3: $\mathcal{A}(\tau)$ is the second from the right side occurrence of $\varnothing$ in $R_{k}$ :


We define $\mathscr{D}$ and $\mathscr{K}$ as in case 1 . The numbers of $p_{*}$ and $\bar{p}_{*}$ in $\mathscr{K}$ are equal. Therefore, the number of $p_{+}$and $\bar{p}_{+}$in $\mathscr{K}$ are also equal. On the other hand, the same is true for $\mathscr{D}$. Contradiction: the number of $\bar{p}_{+}$in $\mathscr{D}$ is the same as in $\mathscr{K}$, but the number of $p_{+}$is greater by one.

Now we construct a proof net $\mathfrak{N}^{\prime}$ for the original sequent $\rightarrow \Gamma: \mathfrak{N}^{\prime} \leftrightharpoons\left\langle\boldsymbol{\Omega}_{\Gamma}, \mathcal{A}^{\prime}, \mathcal{E}^{\prime}\right\rangle$. Here $\mathcal{A}^{\prime}$ contains the arcs of $\mathcal{A}$ that start at old $\otimes$ occurrences, and edges of $\mathcal{E}^{\prime}$ connect those occurrences of $p_{k}$ and $\bar{p}_{k}$ for which the occurrences of $p_{*}$ and $\bar{p}_{*}$ from the corresponding $R_{k}$ and $R_{k}^{\perp}$ are connected by edges of $\mathcal{E}$. It easily follows from the lemmata above that $\mathfrak{N}^{\prime}$ is a proof net for $\rightarrow \Gamma$, therefore MCLL $\vdash \Gamma$. Moreover, if $\mathfrak{N}$ were strong, then $\mathfrak{N}^{\prime}$ is also strong, and $\rightarrow \Gamma$ is derivable in MCLL'.

### 3.3 L $(\backslash ; p)$-grammars and $L^{*}(\backslash ; p)$-grammars

Theorem 24. Any context-free language without the empty word is generated by an $\mathrm{L}\left(\backslash ; p_{1}\right)$-grammar. Any context-free language is generated by an $\mathrm{L}^{*}\left(\backslash ; p_{1}\right)$-grammar.

Proof. Due to Theorems 2 and 15 it is sufficient to show that any $\mathrm{L}(\backslash)$-language (resp., $L^{*}(\backslash)$-language) is generated by an $\mathrm{L}\left(\backslash ; p_{1}\right)$-grammar (resp., $\mathrm{L}^{*}\left(\backslash ; p_{1}\right)$-grammar). Let $M$ be a language generated by an $\mathrm{L}(\backslash)$-grammar (resp., an $\mathrm{L}^{*}(\backslash)$-grammar) $\mathcal{G}$. The grammar is finite; let $N$ be the maximal subscript of a primitive type used in $\mathcal{G}$. Then $\mathcal{G}$ is an $\mathrm{L}\left(\backslash ; p_{1}, \ldots, p_{N}\right)$-grammar (an $\mathrm{L}^{*}\left(\backslash ; p_{1}, \ldots, p_{N}\right)$-grammar). Consider the grammar $\mathcal{G}\left[p_{1}:=\right.$ $\left.A_{1}, \ldots, p_{N}:=A_{N}\right]$. Due to Theorem 20 (resp., Theorem 21) this grammar generates the same language $M$.

Thus, the class of $\mathrm{L}\left(\backslash ; p_{1}\right)$-languages coincides with the class of context-free languages without the empty words and the class of $\mathrm{L}^{*}\left(\backslash ; p_{1}\right)$-languages coincides with the class of all context-free languages.

This construction (or the less complicated Métayer's one) and Theorem 6 yield that any context-free language without the empty word is generated by an $\mathrm{L}\left(\backslash, / ; p_{1}\right)$-grammar with single type assignment.

### 3.4 NP-completeness of Derivability Problems for $\mathrm{L}(\cdot, \backslash ; p), \mathrm{L}(\backslash, / ; p)$, and $\mathrm{L}^{*}(\backslash, / ; p)$

Theorem 25. Derivability problems for $\mathrm{L}(\cdot, \backslash ; p)$, $\mathrm{L}^{*}(\cdot, \backslash ; p), \mathrm{L}(\backslash, / ; p)$, and $\mathrm{L}^{*}(\backslash, / ; p)$ (and, therefore, in their conservative extensions $\mathrm{L}(p), \mathrm{L}^{*}(p), \mathrm{L}_{\mathbf{1}}(\cdot, \backslash ; p), \mathrm{L}_{\mathbf{1}}(\backslash, / ; p)$, and $\left.\mathrm{L}_{1}(p)\right)$ are $N P$-complete.

Proof. Let $\mathscr{L}$ be one of the calculi $\mathrm{L}(\cdot, \backslash), \mathrm{L}^{*}(\cdot, \backslash), \mathrm{L}(\backslash, /)$, or $\mathrm{L}^{*}(\backslash, /)$. We construct a polynomially computable function $\phi$ that reduces derivability in $\mathscr{L}$ to derivability in $\mathscr{L}(p)$. Let $\Pi \rightarrow C$ be an arbitrary sequent and let $N_{\Pi \rightarrow C}$ be the maximal subscript of a primitive type occurring in $\Pi \rightarrow C$. For $N=N_{\Pi \rightarrow C}$ build the types $A_{1}, A_{2}, \ldots, A_{N}$ as shown in Section 3.1. Let $\phi(\Pi \rightarrow C) \leftrightharpoons(\Pi \rightarrow C)\left[p_{1}:=A_{1}, p_{2}:=A_{2}, \ldots, p_{N}:=A_{N}\right]$. Due to conservativity Theorems 22 and 23 yield $\mathscr{L} \vdash \Pi \rightarrow C \Longleftrightarrow \mathscr{L}(p) \vdash \phi(\Pi \rightarrow C)$. Thus, $\phi$ is indeed the reducing function and, since the derivability problem for $\mathscr{L}$ is NP-complete (Theorem 8), so is the derivability problem for $\mathscr{L}(p)$.

For calculi with two divisions one can also use the Métayer construction.
Derivability problems for $\mathrm{L}(\backslash ; p), \mathrm{L}^{*}(\backslash ; p)$, and $\mathrm{L}_{\mathbf{1}}(\backslash ; p)$ are, of course, decidable in polynomial time due to conservativity and Theorems 9 and 19.

### 3.5 The Uniform Substitution

Though for Theorems 24 and 25 it was sufficient to have a faithful substitution acting on a finite set of primitive types, it is interesting whether there exists a uniform substitution, that reduces derivability in $\mathrm{L}(\backslash)$ to derivability in $\mathrm{L}\left(\backslash ; p_{1}\right)$ (resp., $\mathrm{L}^{*}(\backslash)$ to $\left.\mathrm{L}^{*}\left(\backslash ; p_{1}\right)\right)$ independently from the number of primitive types used in the sequent. We build such a substitution using the existing construction. First note that the proof allows parameters to be added to the sequent. Namely, the following statement is true: for any sequent $\Pi \rightarrow C$, where $\Pi=B_{1} \ldots B_{m} \quad B_{1}, \ldots, B_{m}, C \in \operatorname{Tp}\left(\backslash ; p_{1}, \ldots, p_{N}, p_{N+1}, \ldots, p_{N+M}\right)$, we have

$$
\begin{gathered}
\mathrm{L}(\backslash) \vdash \Pi \rightarrow C \Longleftrightarrow \mathrm{~L}(\backslash) \vdash(\Pi \rightarrow C)\left[p_{1}:=A_{1}, \ldots, p_{N}:=A_{N}\right] ; \\
\mathrm{L}^{*}(\backslash) \vdash \Pi \rightarrow C \Longleftrightarrow \mathrm{~L}^{*}(\backslash) \vdash(\Pi \rightarrow C)\left[p_{1}:=A_{1}, \ldots, p_{N}:=A_{N}\right] .
\end{gathered}
$$

Here $p_{N+1}, \ldots, p_{N+M}$ are parameters, that are not affected by the substitution. During the transformation of the (strong) proof net links, connecting their occurrences, are just copied from $\mathfrak{N}$ to $\mathfrak{N}^{\prime}$.

Consider $N=2$. Our construction gives the following types $A_{1}$ and $A_{2}$ (we continue using the notation "multiplication in the denominator"):

$$
\begin{aligned}
& A_{1}=(p \cdot p \cdot(((p \cdot p) \backslash p) \backslash p) \cdot p \cdot p) \backslash p, \\
& A_{2}=(p \cdot p \cdot p \cdot(((p \cdot p) \backslash p) \backslash p) \cdot p) \backslash p .
\end{aligned}
$$

Define two functions from $\operatorname{Tp}(\backslash)$ to $\operatorname{Tp}(\backslash): \mathcal{L}: E \mapsto A_{1}\left[p_{1}:=E\right]$ and $\mathcal{R}: E \mapsto A_{2}\left[p_{1}:=\right.$ $E]$. By renaming primitive types, we get the following:

Lemma 3.8. For any sequent $\Pi \rightarrow C$, where $\Pi=B_{1} \ldots B_{m}$, and $B_{1}, \ldots, B_{m}, C \in$ $\operatorname{Tp}\left(\backslash ; p_{1}, \ldots, p_{n-1}, p_{n}, p_{n+1}\right)$, we have

$$
\begin{aligned}
\mathrm{L}(\backslash) \vdash \Pi \rightarrow C & \Longleftrightarrow \mathrm{~L}(\backslash) \vdash(\Pi \rightarrow C)\left[p_{n}=\mathcal{L}\left(p_{n}\right), p_{n+1}:=\mathcal{R}\left(p_{n}\right)\right] \\
\mathrm{L}^{*}(\backslash) \vdash \Pi \rightarrow C & \Longleftrightarrow \mathrm{~L}^{*}(\backslash) \vdash(\Pi \rightarrow C)\left[p_{n}=\mathcal{L}\left(p_{n}\right), p_{n+1}:=\mathcal{R}\left(p_{n}\right)\right] .
\end{aligned}
$$

(Here $p_{1}, \ldots, p_{n-1}$ act as parameters.)
Now by induction we get the following statement:
Lemma 3.9. For any $n \geqslant 1$ and for any sequent $\Pi \rightarrow C$, where $\Pi=B_{1} \ldots B_{m}$ $B_{1}, \ldots, B_{m}, C \in \operatorname{Tp}\left(\backslash ; p_{1}, \ldots, p_{n}\right)$, the following holds:

$$
\begin{aligned}
& \mathrm{L}(\backslash) \vdash \Pi \rightarrow C \Longleftrightarrow \mathrm{~L}\left(\backslash ; p_{1}\right) \vdash(\Pi \rightarrow C)\left[p_{1}:=\mathcal{L}\left(p_{1}\right),\right. \\
& p_{2}\left.:=\mathcal{L}\left(\mathcal{R}\left(p_{1}\right)\right), \ldots, p_{n-1}:=\mathcal{L}\left(\mathcal{R}^{n-2}\left(p_{1}\right)\right), p_{n}:=\mathcal{R}^{n-1}\left(p_{1}\right)\right] ; \\
& \mathrm{L}^{*}(\backslash) \vdash \Pi \rightarrow C \Longleftrightarrow \mathrm{~L}^{*}\left(\backslash ; p_{1}\right) \vdash(\Pi \rightarrow C)\left[p_{1}:=\mathcal{L}\left(p_{1}\right),\right. \\
& p_{2}\left.:=\mathcal{L}\left(\mathcal{R}\left(p_{1}\right)\right), \ldots, p_{n-1}:=\mathcal{L}\left(\mathcal{R}^{n-2}\left(p_{1}\right)\right), p_{n}:=\mathcal{R}^{n-1}\left(p_{1}\right)\right] .
\end{aligned}
$$

Proof. Proceed by induction on $n$. The induction base is trivial. Suppose our lemma is true for $n$ and prove it for $n+1$. By Lemma 3.8 we have

$$
\mathrm{L}(\backslash) \vdash \Pi \rightarrow C \Longleftrightarrow \mathrm{~L}(\backslash) \vdash(\Pi \rightarrow C)\left[p_{n}:=\mathcal{L}\left(p_{n}\right), p_{n+1}:=\mathcal{R}\left(p_{n}\right)\right] .
$$

The sequent from the right side contains only primitive types $p_{1}, \ldots, p_{n-1}, p_{n}$. Therefore, by induction hypothesis its derivability is equivalent to

$$
\begin{aligned}
\mathrm{L}(\backslash) \vdash(\Pi \rightarrow C)\left[p_{n}:=\mathcal{L}\left(p_{n}\right),\right. & \left.p_{n+1}:=\mathcal{R}\left(p_{n}\right)\right]\left[p_{1}:=\mathcal{L}\left(p_{1}\right)\right. \\
& \left.p_{2}=\mathcal{L}(\mathcal{R}(p)), \ldots, p_{n-1}:=\mathcal{L}\left(\mathcal{R}^{n-2}(p)\right), p_{n}:=\mathcal{R}^{n-1}(p)\right]
\end{aligned}
$$

and the latter sequent is graphically equal to $(\Pi \rightarrow C)\left[p_{1}:=\mathcal{L}(p), p_{2}=\mathcal{L}(\mathcal{R}(p)), \ldots, p_{n-1}:=\right.$ $\left.\mathcal{L}\left(\mathcal{R}^{n-2}(p)\right), p_{n}:=\mathcal{L}\left(\mathcal{R}^{n-1}(p)\right), p_{n+1}:=\mathcal{R}^{n}(p)\right]$.

The $L^{*}(\backslash)$ is handled in the same way.
Theorem 26. For any sequent $\Pi \rightarrow C$, where $\Pi=B_{1} \ldots B_{m}, B_{1}, \ldots, B_{m}, C \in \mathrm{Tp}(\backslash)$, the following equivalences hold:

$$
\begin{aligned}
& \mathrm{L}(\backslash) \vdash \Pi \rightarrow C \Longleftrightarrow \mathrm{~L}\left(\backslash ; p_{1}\right) \vdash(\Pi \rightarrow C) {\left[p_{1}:=\mathcal{L}(p), p_{2}:=\mathcal{L}(\mathcal{R}(p)),\right.} \\
&\left.p_{3}:=\mathcal{L}(\mathcal{R}(\mathcal{R}(p))), \ldots, p_{k}:=\mathcal{L}\left(\mathcal{R}^{k-1}(p)\right), \ldots\right] ; \\
& \mathrm{L}^{*}(\backslash) \vdash \Pi \rightarrow C \Longleftrightarrow \mathrm{~L}^{*}\left(\backslash ; p_{1}\right) \vdash(\Pi \rightarrow C)\left[p_{1}:=\mathcal{L}(p), p_{2}:=\mathcal{L}(\mathcal{R}(p)),\right. \\
&\left.p_{3}:=\mathcal{L}(\mathcal{R}(\mathcal{R}(p))), \ldots, p_{k}:=\mathcal{L}\left(\mathcal{R}^{k-1}(p)\right), \ldots\right] .
\end{aligned}
$$

Proof. Use the previous lemma with $n=N+1$, where $N$ is the greatest subscript of a primitive type occurring in $\Pi \rightarrow C$.

A disadvantage of this construction is the fact that the sequent length grows exponentially: for $p_{k}$ we substitute the type $\mathcal{L}\left(\mathcal{R}^{k-1}(p)\right)$ that contains $9^{k}$ occurrences of $p$. Hence this substitution cannot be used to prove NP-completeness of fragments of the Lambek calculus with one primitive type and at least two connectives.

## Chapter 4

## The Lambek Calculus with the Reversal Operation

### 4.1 L-models for L

In this section and the following three sections the alphabet $\Sigma$ can be finite or countable. Three connectives of the Lambek calculus naturally correspond to three operations on languages without the empty word $\left(M, N \subseteq \Sigma^{+}\right): M \cdot N \leftrightharpoons\{u v \mid u \in M, v \in N\}$, $M \backslash N \leftrightharpoons\left\{u \in \Sigma^{+} \mid(\forall v \in M) v u \in N\right\}, N / M \leftrightharpoons\left\{u \in \Sigma^{+} \mid(\forall v \in M) u v \in N\right\}$.

Definition. An L-model is a pair $\mathcal{M}=\langle\Sigma, w\rangle$, where $\Sigma$ is an alphabet, and $w$ is a mapping of Lambek types into formal languages over $\Sigma$ without the empty word, such that for any $A, B \in \mathrm{Tp}$ the following holds: $w(A \cdot B)=w(A) \cdot w(B), w(A \backslash B)=w(A) \backslash w(B)$ $w(B / A)=w(B) / w(A)$.

This mapping can be defined on primitive types in an arbitrary way, and then it is uniquely propagated to all types.

Since $\Sigma^{+}$with the concatenation operation is the free semigroup generated by $\Sigma$, L-models are also called models on subsets of free semigroups.

Definition. A sequent $F \rightarrow G$ is considered true in a model $\mathcal{M}=\langle\Sigma, w\rangle(\mathcal{M} \vDash F \rightarrow G)$ if $w(F) \subseteq w(G)$.

L-models give sound and complete semantics for L, due to the following theorem:
Theorem 27 (M. Pentus, 1995). A sequent $F \rightarrow G$ is provable in L if and only if it is true in all L-models.

This theorem is proved in [31]; its special case for the product-free fragment (where we keep only types without multiplication) is much easier and appears in [3]. (The notion of truth in an L-model and this theorem can be easily generalized to sequents with more than one type on the left, since $\mathrm{L} \vdash F_{1} F_{2} \ldots F_{n} \rightarrow G$ if and only if $L \vdash F_{1} \cdot F_{2} \cdot \ldots \cdot F_{n} \rightarrow G$.)

### 4.2 The LR Calculus

Now let us consider an extra operation on languages, the reversal. For $u=a_{1} a_{2} \ldots a_{n}$ $\left(a_{1}, \ldots, a_{n} \in \Sigma, n \geq 1\right)$ let $u^{\mathrm{R}} \leftrightharpoons a_{n} \ldots a_{2} a_{1}$, and for $M \subseteq \Sigma^{+}$let $M^{\mathrm{R}} \leftrightharpoons\left\{u^{\mathrm{R}} \mid u \in\right.$ $M\}$. Let us enrich the language of the Lambek calculus with a new unary connective ${ }^{\mathrm{R}}$
(written in the postfix form, $A^{\mathrm{R}}$ ). We shall denote the extended set of types by $\mathrm{Tp}^{\mathrm{R}}$. If $\Gamma=A_{1} A_{2} \ldots A_{n}$, then $\Gamma^{\mathrm{R}} \leftrightharpoons A_{n}^{\mathrm{R}} \ldots A_{2}^{\mathrm{R}} A_{1}^{\mathrm{R}}$.

The notion of L-model is also easily adapted to the new language by adding an additional constraint on $w: w\left(A^{\mathrm{R}}\right)=w(A)^{\mathrm{R}}$.

The calculus $L^{R}$ is obtained from $L$ by adding three new rules for ${ }^{R}$ :

$$
\frac{\Gamma \rightarrow C}{\Gamma^{\mathrm{R}} \rightarrow C^{\mathrm{R}}}\left({ }^{\mathrm{R}} \rightarrow{ }^{\mathrm{R}}\right) \quad \frac{\Gamma A^{\mathrm{RR}} \Delta \rightarrow C}{\Gamma A \Delta \rightarrow C}\left({ }^{\mathrm{RR}} \rightarrow\right)_{\mathrm{E}} \quad \frac{\Gamma \rightarrow C^{\mathrm{RR}}}{\Gamma \rightarrow C}\left(\rightarrow^{\mathrm{RR}}\right)_{\mathrm{E}}
$$

It is easy to see that $\mathrm{L}^{\mathrm{R}}$ is sound with respect to L -models.
Proposition 4.1. The calculus $\mathrm{L}^{\mathrm{R}}$ is a conservative extension of L .
Proof. The "if" part is obvious. The "only if" part follows from L-completeness of L and L-soundness of $\mathrm{L}^{\mathrm{R}}$ : if $F \rightarrow G$ is provable in $\mathrm{L}^{\mathrm{R}}$, then it is true in all L-models, and, therefore, is provable in L .

L-completeness for the product-free fragment of $\mathrm{L}^{\mathrm{R}}$ is proved in [29] by a modification of Buszkowski's argument [3] (in [29] the reversal connective is called involution and denoted by ${ }^{`}$ instead of ${ }^{R}$; the calculus is formulated in a different, but equivalent way). In [29] one can also find a proof of L-completeness of the division-free fragment (where only • and ${ }^{\mathrm{R}}$ connectives are kept). We shall prove L-completeness of the whole calculus.

Theorem 28. A sequent $F \rightarrow G\left(F, G \in \mathrm{Tp}^{\mathrm{R}}\right)$ is provable in $\mathrm{L}^{\mathrm{R}}$ if and only if it is true in all L-models.

A variant of this calculus that allows empty antecedents (an extension with the ${ }^{R}$ connective of $L^{*}$, the variant of $L$ without the restriction $\Pi \neq \Lambda$ on the $(\rightarrow \backslash)$ and $(\rightarrow /)$ rules) is presented in [17]. The calculus $L^{*}$ itself is complete with respect to L-models allowing empty words in the languages (free monoid models) [24], but L-completeness of its extension with the ${ }^{\mathrm{R}}$ connective is still an open problem.

### 4.3 Normal Form for $L^{\mathrm{R}}$ types

The relation $\leftrightarrow_{\mathrm{L}^{\mathrm{R}}}$ is an equivalence relation. Moreover, it is a congruence relation with respect to all connectives (due to the $\left({ }^{\mathrm{R}} \rightarrow^{\mathrm{R}}\right)$ rule for ${ }^{\mathrm{R}}$ and Proposition 1.2 for others). The following lemma is checked explicitly by presenting the corresponding derivations in $\mathrm{L}^{\mathrm{R}}$ :

Lemma 4.2. The following equivalences hold in $\mathrm{L}^{\mathrm{R}}$ :

1. $(A \cdot B)^{\mathrm{R}} \leftrightarrow B^{\mathrm{R}} \cdot A^{\mathrm{R}}$;
2. $(A \backslash B)^{\mathrm{R}} \leftrightarrow B^{\mathrm{R}} / A^{\mathrm{R}}$;
3. $(B / A)^{\mathrm{R}} \leftrightarrow A^{\mathrm{R}} \backslash B^{\mathrm{R}}$;
4. $A^{\mathrm{RR}} \leftrightarrow A$.

Definition. For $A \in \mathrm{Tp}^{\mathrm{R}}$ we define $\operatorname{tr}(A)$ by induction on the number of connectives in A:

1. $\operatorname{tr}\left(p_{i}\right) \leftrightharpoons p_{i}$;
2. $\operatorname{tr}\left(p_{i}^{\mathrm{R}}\right) \leftrightharpoons p_{i}^{\mathrm{R}}$;
3. $\operatorname{tr}(A \cdot B) \leftrightharpoons \operatorname{tr}(A) \cdot \operatorname{tr}(B)$;
4. $\operatorname{tr}(A \backslash B) \leftrightharpoons \operatorname{tr}(A) \backslash \operatorname{tr}(B)$;
5. $\operatorname{tr}(B / A) \leftrightharpoons \operatorname{tr}(B) / \operatorname{tr}(A)$;
6. $\operatorname{tr}\left((A \cdot B)^{\mathrm{R}}\right) \leftrightharpoons \operatorname{tr}\left(B^{\mathrm{R}}\right) \cdot \operatorname{tr}\left(A^{\mathrm{R}}\right) ;$
7. $\operatorname{tr}\left((A \backslash B)^{\mathrm{R}}\right) \leftrightharpoons \operatorname{tr}\left(B^{\mathrm{R}}\right) / \operatorname{tr}\left(A^{\mathrm{R}}\right)$;
8. $\operatorname{tr}\left((B / A)^{\mathrm{R}}\right) \leftrightharpoons \operatorname{tr}\left(A^{\mathrm{R}}\right) \backslash \operatorname{tr}\left(B^{\mathrm{R}}\right) ;$
9. $\operatorname{tr}\left(A^{\mathrm{RR}}\right) \leftrightharpoons \operatorname{tr}(A)$.

The following statement is proved by induction using Lemma 4.2:
Proposition 4.3. Any $A \in \mathrm{Tp}^{\mathrm{R}}$ is equivalent to $\operatorname{tr}(A)$.
We call $\operatorname{tr}(A)$ the normal form of $A$. In the normal form, the ${ }^{\mathrm{R}}$ connective can appear only on occurrences of primitive types.

### 4.4 L-completeness of $\mathrm{L}^{\mathrm{R}}$ (Proof)

Now we are going to prove Theorem 28 (the "if" part) by contraposition. Let $\mathrm{L}^{\mathrm{R}} \nvdash F_{0} \rightarrow$ $G_{0}$. We need to construct a countermodel for $F_{0} \rightarrow G_{0}$, i.e., a model in which this sequent is not true.

Let $\operatorname{Pr}^{\prime} \leftrightharpoons \operatorname{Pr} \cup\left\{p^{\mathrm{R}} \mid p \in \operatorname{Pr}\right\}$, and let $\mathrm{L}^{\prime}$ be the Lambek calculus with $\operatorname{Pr}^{\prime}$ taken as the set of primitive types instead of $\operatorname{Pr}$. Here ${ }^{\mathrm{R}}$ is not a connective, and $p^{\mathrm{R}}$ is considered just a new primitive type, independent from $p$. Obviously, if $L^{\prime} \vdash F \rightarrow G$, then $\mathrm{L}^{\mathrm{R}} \vdash$ $F \rightarrow G$.

Let $F \leftrightharpoons \operatorname{tr}\left(F_{0}\right), G \leftrightharpoons \operatorname{tr}\left(G_{0}\right)$. Then $\mathrm{L}^{\mathrm{R}} \nvdash F \rightarrow G$, whence $\mathrm{L}^{\prime} \nvdash F \rightarrow G$. The calculus $\mathrm{L}^{\prime}$ is essentially the same as L , therefore Theorem 27 gives us a structure $\mathcal{M}=$ $\langle\Sigma, w\rangle$ such that $w(F) \nsubseteq w(G)$. The structure $\mathcal{M}$ indeed falsifies $F \rightarrow G$, but it is not a model in the sense of our new language: some of the conditions $w\left(p_{i}^{\mathrm{R}}\right)=w\left(p_{i}\right)^{\mathrm{R}}$ might be not satisfied.

Let $\Phi$ be the set of all subtypes of $F \rightarrow G$ (including $F$ and $G$ themselves; the notion of subtype is understood in the sense of $\mathrm{L}^{\mathrm{R}}$ ). The construction of $\mathcal{M}$ (see [31]) guarantees that $w(A) \neq \varnothing$ for all $A \in \Phi$. This is the only specific property of $\mathcal{M}$ we shall need.

We introduce an inductively defined counter $f(A), A \in \Phi: f\left(p_{i}\right) \leftrightharpoons 1, f\left(p_{i}^{\mathrm{R}}\right) \leftrightharpoons 1$, $f(A \cdot B) \leftrightharpoons f(A)+f(B)+10, f(A \backslash B) \leftrightharpoons f(B), f(B / A) \leftrightharpoons f(B)$. Let $K \leftrightharpoons \max \{f(A) \mid$ $A \in \Phi\}, N \leftrightharpoons 2 K+25$ ( $N$ should be odd, greater then $K$, and big enough itself).

Let $\Sigma_{1} \leftrightharpoons \Sigma \times\{1, \ldots, N\}$. We shall denote the pair $\langle a, j\rangle \in \Sigma_{1}$ by $a^{(j)}$. Elements of $\Sigma$ and $\Sigma_{1}$ will be called letters and symbols respectively. A symbol can be even or odd depending on the parity of the superscript. Consider a homomorphism $h: \Sigma^{+} \rightarrow \Sigma_{1}^{+}$, defined as follows: $h(a) \leftrightharpoons a^{(1)} a^{(2)} \ldots a^{(N)}(a \in \Sigma), h\left(a_{1} \ldots a_{n}\right) \leftrightharpoons h\left(a_{1}\right) \ldots h\left(a_{n}\right)$. Let $P \leftrightharpoons h\left(\Sigma^{+}\right)=\left\{a_{1}^{(1)} \ldots a_{1}^{(N)} \ldots a_{n}^{(1)} \ldots a_{n}^{(N)} \mid n \geq 1, a_{i} \in \Sigma\right\}$. Note that $h$ is a bijection between $\Sigma^{+}$and $P$.

Lemma 4.4. For all $M, N \subseteq \Sigma^{+}$we have

1. $h(M \cdot N)=h(M) \cdot h(N)$;
2. if $M \neq \varnothing$, then $h(M \backslash N)=h(M) \backslash h(N)$ and $h(N / M)=h(N) / h(M)$.

Proof.

1. By the definition of a homomorphism.
2. 

$\subseteq$ Let $u \in h(M \backslash N)$. Then $u=h\left(u^{\prime}\right)$ for some $u^{\prime} \in M \backslash N$. For all $v^{\prime} \in M$ we have $v^{\prime} u^{\prime} \in N$. Take an arbitrary $v \in h(M), v=h\left(v^{\prime}\right)$ for some $v^{\prime} \in M$. Since $u^{\prime} \in M \backslash N, v^{\prime} u^{\prime} \in N$, whence $v u=h\left(v^{\prime}\right) h\left(u^{\prime}\right)=h\left(v^{\prime} u^{\prime}\right) \in h(N)$. Therefore $u \in h(M) \backslash h(N)$.
$\supseteq$ Let $u \in h(M) \backslash h(N)$. First we claim that $u \in P$. Suppose the contrary: $u \notin P$. Take $v^{\prime} \in M$ ( $M$ is nonempty by assumption). Since $v=h\left(v^{\prime}\right) \in P$, $v u \notin P$. On the other hand, $v u \in h(N) \subseteq P$. Contradiction. Now, since $u \in P, u=h\left(u^{\prime}\right)$ for some $u^{\prime} \in \Sigma^{+}$. For an arbitrary $v^{\prime} \in M$ and $v \leftrightharpoons h\left(v^{\prime}\right)$ we have $h\left(v^{\prime} u^{\prime}\right)=v u \in h(N)$, whence $v^{\prime} u^{\prime} \in N$, whence $u^{\prime} \in M \backslash N$. Therefore, $u=h\left(u^{\prime}\right) \in h(M \backslash N)$.
The / case is handled symmetrically.

We construct a new model $\mathcal{M}_{1}=\left\langle\Sigma_{1}, w_{1}\right\rangle$, where $w_{1}(z) \leftrightharpoons h(w(z))\left(z \in \operatorname{Pr}^{\prime}\right)$. Due to Lemma 4.4, $w_{1}(A)=h\left(w_{1}(A)\right)$ for all $A \in \Phi$, whence $w_{1}(F)=h(w(F)) \nsubseteq h(w(G))=$ $w_{1}(G)\left(\mathcal{M}_{1}\right.$ is also a countermodel in the language without $\left.{ }^{\mathrm{R}}\right)$.

Now we introduce several auxiliary subsets of $\Sigma_{1}^{+}(\operatorname{byb} \operatorname{Subw}(M)$ we denote the set of all nonempty subwords of words from $M$, i.e. $\operatorname{Subw}(M) \leftrightharpoons\left\{u \in \Sigma_{1}^{+} \mid\left(\exists v_{1}, v_{2} \in\right.\right.$ $\left.\left.\Sigma_{1}^{*}\right) v_{1} u v_{2} \in M\right\}$ ):
$T_{1} \leftrightharpoons\left\{u \in \Sigma_{1}^{+} \mid u \notin \operatorname{Subw}\left(P \cup P^{\mathrm{R}}\right)\right\} ;$
$T_{2} \leftrightharpoons\left\{u \in \operatorname{Subw}\left(P \cup P^{\mathrm{R}}\right) \mid\right.$ the first or the last symbol of $u$ is even $\} ;$
$E \leftrightharpoons\left\{u \in \operatorname{Subw}\left(P \cup P^{\mathrm{R}}\right)-\left(P \cup P^{\mathrm{R}}\right) \mid\right.$ both the first symbol and the last symbol of $u$ are odd\}.

The sets $P, P^{\mathrm{R}}, T_{1}, T_{2}$, and $E$ form a partition of $\Sigma_{1}^{+}$into nonintersecting parts. For example, $a^{(1)} b^{(10)} a^{(2)} \in T_{1}, a^{(N)} b^{(1)} \ldots b^{(N-1)} \in T_{2}, a^{(7)} a^{(6)} a^{(5)} \in E(a, b \in \Sigma)$.

Let $T \leftrightharpoons T_{1} \cup T_{2}, T_{i}(k) \leftrightharpoons\left\{u \in T_{i}| | u \mid \geq k\right\}(i=1,2,|u|$ is the length of $u)$, $T(k) \leftrightharpoons T_{1}(k) \cup T_{2}(k)=\{u \in T| | u \mid \geq k\}$.

Note that if the first or the last symbol of $u$ is even, then it belongs to $T$, no matter whether it belongs to $\operatorname{Subw}\left(P \cup P^{\mathrm{R}}\right)$.

The index $k$ (possibly with subscripts) here and further ranges from 1 to $K$. For all $k$ we have $T(k) \supseteq T(K)$.

## Lemma 4.5.

1. $P \cdot P \subseteq P, P^{\mathrm{R}} \cdot P^{\mathrm{R}} \subseteq P^{\mathrm{R}}$;
2. $T^{\mathrm{R}}=T, T(k)^{\mathrm{R}}=T(k)$;
3. $P \cdot P^{\mathrm{R}} \subseteq T(K), P^{\mathrm{R}} \cdot P \subseteq T(K)$;
4. $P \cdot T \subseteq T(K), T \cdot P \subseteq T(K)$;
5. $P^{\mathrm{R}} \cdot T \subseteq T(K), T \cdot P^{\mathrm{R}} \subseteq T(K)$;
6. $T \cdot T \subseteq T$;

Proof.

1. Obvious.
2. Directly follows from our definitions.
3. Any element of $P \cdot P^{\mathrm{R}}$ or $P^{\mathrm{R}} \cdot P$ does not belong to $\operatorname{Subw}\left(P \cup P^{\mathrm{R}}\right)$ and its length is at least $2 N>K$. Therefore it belongs to $T_{1}(K) \subseteq T(K)$.
4. Let $u \in P$ and $v \in T$. If $v \in T_{1}$, then $u v$ is also in $T_{1}$. Let $v \in T_{2}$. If the last symbol of $v$ is even, then $u v \in T$. If the last symbol of $v$ is odd, then $u v \notin \operatorname{Subw}\left(P \cup P^{\mathrm{R}}\right)$, whence $u v \in T_{1} \subseteq T$. Since $|u v|>|u| \geq N>K, u v \in T(K)$.
The claim $T \cdot P \subseteq T$ is handled symmetrically.
5. $P^{\mathrm{R}} \cdot T=P^{\mathrm{R}} \cdot T^{\mathrm{R}}=(T \cdot P)^{\mathrm{R}} \subseteq T(K)^{\mathrm{R}}=T(K) . T \cdot P^{\mathrm{R}}=T^{\mathrm{R}} \cdot P^{\mathrm{R}}=(P \cdot T)^{\mathrm{R}} \subseteq$ $T(K)^{\mathrm{R}}=T(K)$.
6. Let $u, v \in T$. If at least one of these two words belongs to $T_{1}$, then $u v \in T_{1}$. Let $u, v \in T_{2}$. If the first symbol of $u$ or the last symbol of $v$ is even, then $u v \in T$. In the other case $u$ ends with an even symbol, and $v$ starts with an even symbol. But then we have two consecutive even symbols in $u v$, therefore $u v \in T_{1}$.

Let us call words of the form $a^{(i)} a^{(i+1)} a^{(i+2)}, a^{(N-1)} a^{(N)} b^{(1)}$, and $a^{(N)} b^{(1)} b^{(2)}(a, b \in$ $\Sigma, 1 \leq i \leq N-2$ ) valid triples of type $I$ and their reversals (namely, $a^{(i+2)} a^{(i+1)} a^{(i)}$, $b^{(1)} a^{(N)} a^{(N-1)}$, and $\left.b^{(2)} b^{(1)} a^{(N)}\right)$ valid triples of type II. Note that valid triples of type I (resp., of type II) are the only possible three-symbol subwords of words from $P$ (resp., $P^{\mathrm{R}}$ ).

Lemma 4.6. A word $u$ of length at least three is a subword of a word from $P \cup P^{\mathrm{R}}$ if and only if any three-symbol subword of $u$ is a valid triple of type I or II.

Proof. The nontrivial part is "if". We proceed by induction on $|u|$. Induction base $(|u|=3)$ is trivial. Let $u$ be a word of length $m+1$ satisfying the condition and let $u=u^{\prime} x\left(x \in \Sigma_{1}\right)$. By induction hypothesis $\left(\left|u^{\prime}\right|=m\right), u^{\prime} \in \operatorname{Subw}\left(P \cup P^{\mathrm{R}}\right)$. Let $u^{\prime} \in \operatorname{Subw}(P)$ (the other case is handled symmetrically); $u^{\prime}$ is a subword of some word $v \in P$. Consider the last three symbols of $u$. Since the first two of them also belong to $u^{\prime}$, this three-symbol word is a valid triple of type I, not type II. If it is of the form $a^{(i)} a^{(i+1)} a^{(i+2)}$ or $a^{(N)} b^{(1)} b^{(2)}$, then $x$ coincides with the symbol next to the occurrence of $u^{\prime}$ in $v$, and therefore $u=u^{\prime} x$ is also a subword of $v$. If it is of the form $a^{(N-1)} a^{(N)} b^{(1)}$, then, provided $v=v_{1} u^{\prime} v_{2}, v_{1} u^{\prime}$ is also an element of $P$, and so is the word $v_{1} u^{\prime} b^{(1)} b^{(2)} \ldots b^{(N)}$, which contains $u=u^{\prime} b^{(1)}$ as a subword. Thus, in all cases $u \in \operatorname{Subw}(P)$.

Now we construct one more model $\mathcal{M}_{2}=\left\langle\Sigma_{1}, w_{2}\right\rangle$, where $w_{2}\left(p_{i}\right) \leftrightharpoons w_{1}\left(p_{i}\right) \cup$ $w_{1}\left(p_{i}^{\mathrm{R}}\right)^{\mathrm{R}} \cup T, w_{2}\left(p_{i}^{\mathrm{R}}\right) \leftrightharpoons w_{1}\left(p_{i}\right)^{\mathrm{R}} \cup w_{1}\left(p_{i}^{\mathrm{R}}\right) \cup T$. This model is a model even in the sense of the enriched language. To finish the proof, we need to check that $\mathcal{M}_{2} \not \nexists F \rightarrow G$.

Lemma 4.7. For any $A \in \Phi$ the following holds:

1. $w_{2}(A) \subseteq P \cup P^{\mathrm{R}} \cup T$;
2. $w_{2}(A) \supseteq T(f(A))$;
3. $w_{2}(A) \cap P=w_{1}(A)$ (in particular, $w_{2}(A) \cap P \neq \varnothing$ );
4. $w_{2}(A) \cap P^{\mathrm{R}}=w_{1}\left(\operatorname{tr}\left(A^{\mathrm{R}}\right)\right)^{\mathrm{R}}$ (in particular, $w_{2}(A) \cap P^{\mathrm{R}} \neq \varnothing$ ).

Proof. We prove all the statements simultaneously by induction on type $A$. The induction base is trivial. Further we shall refer to the $i$-th statement of the induction hypothesis ( $i=1,2,3,4)$ as "IH- $i$ ".

1. Consider three possible cases.
a) $A=B \cdot C$. Then $w_{2}(A)=w_{2}(B) \cdot w_{2}(C) \subseteq\left(P \cup P^{\mathrm{R}} \cup T\right) \cdot\left(P \cup P^{\mathrm{R}} \cup T\right) \subseteq P \cup P^{\mathrm{R}} \cup T$ (Lemma 4.5).
b) $A=B \backslash C$. Suppose the contrary: in $w_{2}(A)$ there exists an element $u \in E$. Then $v u \in w_{2}(C)$ for any $v \in w_{2}(B)$. We consider several subcases and show that each of those leads to a contradiction.
i) $u \in \operatorname{Subw}(P)$, and the superscript of the first symbol of $u$ is not 1 . Let the first symbol of $u$ be $a^{(i)}$. Note that $i$ is odd and $i>2$. Take $v=a^{(3)} \ldots a^{(N)} a^{(1)} \ldots a^{(i-1)}$. The word $v$ has length at least $N \geq K$ and ends with an even symbol, therefore $v \in T(K) \subseteq$ $T(f(B)) \subseteq w_{2}(B)(\mathrm{IH}-2)$. On the other hand, $v u \in \operatorname{Subw}(P)$ and the first symbol and the last symbol of $v u$ are odd. Therefore, $v u \in E$ and $v u \in w_{2}(C)$, but $w_{2}(C) \cap E=\varnothing$ (IH-1). Contradiction.
ii) $u \in \operatorname{Subw}(P)$, and the first symbol of $u$ is $a^{(1)}$ (then the superscript of the last symbol of $u$ is not $N$, because otherwise $u \in P$ ). Take $v \in w_{2}(B) \cap P$ (this set is nonempty due to IH-3). The first and the last symbol of $v u$ is odd, and $v u \in \operatorname{Subw}(P)-P$, therefore $v u \in E$. Contradiction.
iii) $u \in \operatorname{Subw}\left(P^{\mathrm{R}}\right)$, and the superscript of the first symbol of $u$ is not $N$ (the first symbol of $u$ is $a^{(i)}, i$ is odd). Take $v=a^{(N-2)} \ldots a^{(1)} a^{(N)} \ldots a^{(i+1)} \in T(K) \subseteq w_{2}(B)$. Again, $v u \in E$.
iv) $u \in \operatorname{Subw}\left(P^{\mathrm{R}}\right)$, and the first symbol of $u$ is $a^{(N)}$. Take $v \in w_{2}(B) \cap P^{\mathrm{R}}$ (nonempty due to IH-4). $v u \in E$.
c) $A=C / B$. Proceed symmetrically.
2. Consider three possible cases.
a) $A=B \cdot C$. Let $k_{1} \leftrightharpoons f(B), k_{2} \leftrightharpoons f(C), k \leftrightharpoons k_{1}+k_{2}+10=f(A)$. Due to IH-2, $w_{2}(B) \supseteq T\left(k_{1}\right)$ and $w_{2}(C) \supseteq T\left(k_{2}\right)$. Take $u \in T(k)$. We have to prove that $u \in w_{2}(A)$. Consider several subcases.
i) $u \in T_{1}(k)$. By Lemma $4.6\left(|u| \geq k>3\right.$ and $\left.u \notin \operatorname{Subw}\left(P \cup P^{\mathrm{R}}\right)\right)$ in $u$ there is a three-symbol subword $x y z$ that is not a valid triple of type I or II. Divide the word $u$ into two parts, $u=u_{1} u_{2}$, such that $\left|u_{1}\right| \geq k_{1}+5,\left|u_{2}\right| \geq k_{2}+5$. If needed, shift the border between parts by one symbol to the left or to the right, so that the subword $x y z$ lies entirely in one part. Let this part be $u_{2}$ (the other case is handled symmetrically). Then $u_{2} \in T_{1}\left(k_{2}\right)$. If $u_{1}$ is also in $T_{1}$, then the proof is finished. Consider the other case. Note that in any word from $\operatorname{Subw}\left(P \cup P^{\mathrm{R}}\right)$ among any three consecutive symbols at least one is even. Shift the border to the left by at most 2 symbols to make the last symbol of $u_{1}$ even. Then $u_{1} \in T\left(k_{1}\right)$, and $u_{2}$ remains in $T_{1}\left(k_{2}\right)$. Thus $u=u_{1} u_{2} \in T\left(k_{1}\right) \cdot T\left(k_{2}\right) \subseteq w_{2}(B) \cdot w_{2}(C)=w_{2}(A)$.
ii) $u \in T_{2}(k)$. Let $u$ end with an even symbol (the other case is symmetric). Divide the word $u$ into two parts, $u=u_{1} u_{2},\left|u_{1}\right| \geq k_{1}+5, u_{2} \geq k_{2}+5$, and shift the border (if needed), so that the last symbol of $u_{1}$ is even. Then both $u_{1}$ and $u_{2}$ end with an even symbol, and therefore $u_{1} \in T\left(k_{1}\right)$ and $u_{2} \in T\left(k_{2}\right)$.
b) $A=B \backslash C$. Let $k \leftrightharpoons f(C)=f(A)$. By IH-2, $w_{2}(C) \supseteq T(k)$. Take $u \in T(k)$ and an arbitrary $v \in w_{2}(B) \subseteq P \cup P^{\mathrm{R}} \cup T$. By Lemma 4.5, statements 4-6, vu $\in$ $\left(P \cup P^{\mathrm{R}} \cup T\right) \cdot T \subseteq T$, and since $|v u|>|u| \geq k, v u \in T(k) \subseteq w_{2}(C)$. Thus $u \in w_{2}(A)$.
c) $A=C / B$. Symmetrically.
3. Consider three possible cases.
a) $A=B \cdot C$.
$\supseteq u \in w_{1}(A)=w_{1}(B) \cdot w_{1}(C) \subseteq w_{2}(B) \cdot w_{2}(C)=w_{2}(A)$ (IH-3); $u \in P$.
$\subseteq$ Suppose $u \in P$ and $u \in w_{2}(A)=w_{2}(B) \cdot w_{2}(C)$. Then $u=u_{1} u_{2}$, where $u_{1} \in$ $w_{2}(B)$ and $u_{2} \in w_{2}(C)$. First we claim that $u_{1} \in P$. Suppose the contrary, $u_{1} \notin P$. By IH$1, u_{1} \in P^{\mathrm{R}} \cup T, u_{2} \in P \cup P^{\mathrm{R}} \cup T$, and therefore $u=u_{1} u_{2} \in\left(P^{\mathrm{R}} \cup T\right) \cdot\left(P \cup P^{\mathrm{R}} \cup T\right) \subseteq P^{\mathrm{R}} \cup T$ (Lemma 4.5, statements 1, 3-6). Hence $u \notin P$. Contradiction. Thus, $u_{1} \in P$. Similarly, $u_{2} \in P$, and by IH-3 we obtain $u_{1} \in w_{1}(B)$ and $u_{2} \in w_{1}(C)$, whence $u=u_{1} u_{2} \in w_{1}(A)$.
b) $A=B \backslash C$.
$\supseteq$ Take $u \in w_{1}(B \backslash C)$. For any $v \in w_{1}(B)$ we have $v u \in w_{1}(C)$. We claim that $u \in w_{2}(B \backslash C)$. Take $v \in w_{2}(B) \subseteq P \cup P^{\mathrm{R}} \cup T$ (IH-1). If $v \in P$, then $v \in w_{1}(B)(\mathrm{IH}-3)$, and $v u \in w_{1}(C) \subseteq w_{2}(C)(\mathrm{IH}-3)$. If $v \in P^{\mathrm{R}} \cup T$, then $v u \in\left(P^{\mathrm{R}} \cup T\right) \cdot P \subseteq T(K) \subseteq w_{2}(C)$ (Lemma 4.5, statements 3 and 4, and IH-2). Therefore, $u \in w_{2}(B) \backslash w_{2}(C)=w_{2}(B \backslash C)$; also we have $u \in P$, since $w_{1}(B \backslash C) \subseteq P$.
$\subseteq$ If $u \in w_{2}(B \backslash C)$ and $u \in P$, then for any $v \in w_{1}(B) \subseteq w_{2}(B)$ we have $v u \in w_{2}(C)$. Since $v, u \in P, v u \in P$. By IH-3, $v u \in w_{1}(C)$. Thus $u \in w_{1}(B \backslash C)$.
c) $A=C / B$. Symmetrically.
4. Consider three cases.
a) $A=B \cdot C$. Then $\operatorname{tr}\left(A^{\mathrm{R}}\right)=\operatorname{tr}\left(C^{\mathrm{R}}\right) \cdot \operatorname{tr}\left(B^{\mathrm{R}}\right)$.
$\supseteq u \in w_{1}\left(\operatorname{tr}\left(A^{\mathrm{R}}\right)\right)^{\mathrm{R}}=w_{1}\left(\operatorname{tr}\left(C^{\mathrm{R}}\right) \cdot \operatorname{tr}\left(B^{\mathrm{R}}\right)\right)^{\mathrm{R}}=\left(w_{1}\left(\operatorname{tr}\left(C^{R}\right)\right) \cdot w_{1}\left(\operatorname{tr}\left(B^{\mathrm{R}}\right)\right)\right)^{\mathrm{R}}=$ $w_{1}\left(\operatorname{tr}\left(B^{\mathrm{R}}\right)\right)^{\mathrm{R}} \cdot w_{1}\left(\operatorname{tr}\left(C^{\mathrm{R}}\right)\right)^{\mathrm{R}} \subseteq w_{2}(B) \cdot w_{2}(C)=w_{2}(A)(\mathrm{IH}-4) ; u \in P^{\mathrm{R}}$.
$\subseteq$ Let $u \in P^{\mathrm{R}}$ and $u \in w_{2}(A)=w_{2}(B) \cdot w_{2}(C)$. Then $u=u_{1} u_{2}$, where $u_{1} \in w_{2}(B)$, $u_{2} \in w_{2}(C)$. We claim that $u_{1}, u_{2} \in P^{\mathrm{R}}$. Suppose the contrary: $u_{1} \notin P^{\mathrm{R}}$. Then $u_{1} \in P \cup T(\mathrm{IH}-1), u_{2} \in P \cup P^{\mathrm{R}} \cup T$, whence $u=u_{1} u_{2} \in(P \cup T) \cdot\left(P \cup P^{\mathrm{R}} \cup T\right) \subseteq P \cup T$. Contradiction $\left(u \in P^{\mathrm{R}}\right)$. Thus, $u_{1} \in P^{\mathrm{R}}$, and therefore $u_{2} \in P^{R}$, and, using IH-4, we obtain $u_{1} \in w_{1}\left(\operatorname{tr}\left(B^{\mathrm{R}}\right)\right)^{\mathrm{R}}, u_{2} \in w_{1}\left(\operatorname{tr}\left(C^{\mathrm{R}}\right)\right)^{\mathrm{R}}$. Hence $u=u_{1} u_{2} \in w_{1}\left(\operatorname{tr}\left(B^{\mathrm{R}}\right)\right)^{\mathrm{R}}$. $w_{1}\left(\operatorname{tr}\left(C^{\mathrm{R}}\right)\right)^{\mathrm{R}}=\left(w_{1}\left(\operatorname{tr}\left(C^{\mathrm{R}}\right)\right) \cdot w_{1}\left(\operatorname{tr}\left(B^{\mathrm{R}}\right)\right)\right)^{\mathrm{R}}=w_{1}\left(\operatorname{tr}\left(C^{\mathrm{R}}\right) \cdot \operatorname{tr}\left(B^{\mathrm{R}}\right)\right)^{\mathrm{R}}=w_{1}\left(\operatorname{tr}\left(A^{\mathrm{R}}\right)\right)^{\mathrm{R}}$.
b) $A=B \backslash C$. Then $\operatorname{tr}\left(A^{\mathrm{R}}\right)=\operatorname{tr}\left(C^{\mathrm{R}}\right) / \operatorname{tr}\left(B^{\mathrm{R}}\right)$.
$\supseteq$ Let $u \in w_{1}\left(\operatorname{tr}\left(C^{\mathrm{R}}\right) / \operatorname{tr}\left(B^{\mathrm{R}}\right)\right)^{\mathrm{R}}=w_{1}\left(\operatorname{tr}\left(B^{\mathrm{R}}\right)\right)^{\mathrm{R}} \backslash w_{1}\left(\operatorname{tr}\left(C^{\mathrm{R}}\right)\right)^{\mathrm{R}}$, so for every $v \in$ $w_{1}\left(\operatorname{tr}\left(B^{\mathrm{R}}\right)\right)^{\mathrm{R}}$ we have $v u \in w_{1}\left(\operatorname{tr}\left(C^{\mathrm{R}}\right)\right)^{\mathrm{R}}$. We claim that $u \in w_{2}(B \backslash C)$. Take an arbitrary $v \in w_{2}(B) \subseteq P \cup P^{\mathrm{R}} \cup T$ (IH-1). If $v \in P^{\mathrm{R}}$, then $v \in w_{1}\left(\operatorname{tr}\left(B^{\mathrm{R}}\right)\right)^{\mathrm{R}}$ (IH-4), whence $v u \in w_{1}\left(\operatorname{tr}\left(C^{\mathrm{R}}\right)\right)^{\mathrm{R}} \subseteq w_{2}(C)$.
If $v \in P \cup T$, then (since $u \in P^{\mathrm{R}}$ ) we have $v u \in(P \cup T) \cdot P^{\mathrm{R}} \subseteq T(K) \subseteq w_{2}(C)$ (Lemma 4.5 and IH-2).
$\subseteq$ If $u \in w_{2}(B \backslash C)$ and $u \in P^{\mathrm{R}}$, then for any $v \in w_{1}\left(\operatorname{tr}\left(B^{\mathrm{R}}\right)\right)^{\mathrm{R}} \subseteq w_{2}(B)$ we have $v u \in w_{2}(C)$. Since $v, u \in P^{R}, v u \in P^{\mathrm{R}}$, therefore $v u \in w_{1}\left(\operatorname{tr}\left(C^{\mathrm{R}}\right)\right)^{\mathrm{R}}(\mathrm{IH}-4)$. Thus $u \in w_{1}\left(\operatorname{tr}\left(B^{\mathrm{R}}\right)\right)^{\mathrm{R}} \backslash w_{1}\left(\operatorname{tr}\left(C^{\mathrm{R}}\right)\right)^{\mathrm{R}}=w_{1}\left(A^{\mathrm{R}}\right)^{\mathrm{R}}$.
c) $A=C / B$. Symmetrically.

This completes the proof of Lemma 4.7.
Since $w_{1}(F) \nsubseteq w_{1}(G)$, there exists an element $u_{0}$ such that $u_{0} \in w_{1}(F)$ and $u_{0} \notin$ $w_{1}(G)$. Since $u_{0} \in P, u_{0} \in w_{2}(F)$ and $u_{0} \notin w_{2}(G)$. Therefore, $w_{2}(F) \nsubseteq w_{2}(G)$. Since $F_{0} \leftrightarrow F, G_{0} \leftrightarrow G$, and $\mathrm{L}^{\mathrm{R}}$ is L-sound, we see that $w_{2}\left(F_{0}\right)=w_{2}(F), w_{2}\left(G_{0}\right)=w_{2}(G)$, and $\mathcal{M}_{2}$ is a countermodel for $F_{0} \rightarrow G_{0}$. This completes the proof of Theorem 28.

Note that we have constructed a countermodel (in the sense of the extended language) for any sequent $F \rightarrow G$ that is not provable in $\mathrm{L}^{\prime}$ (this could be potentially weaker than $\left.\mathrm{L}^{\mathrm{R}} \nvdash F \rightarrow G\right)$. Thus we get the following statement:

Proposition 4.8. $\mathrm{L}^{\mathrm{R}} \vdash A_{1} \ldots A_{n} \rightarrow B$ if and only if $\mathrm{L}^{\prime} \vdash \operatorname{tr}\left(A_{1}\right) \ldots \operatorname{tr}\left(A_{n}\right) \rightarrow \operatorname{tr}(B)$.

## $4.5 \quad L^{R}$ : Grammars and Complexity

Theorem 29. The class of all $\mathrm{L}^{\mathrm{R}}$-languages coincides with the class of all context-free languages without the empty word.

Proof. Every context-free language is an $\mathrm{L}^{\mathrm{R}}$-language due to Theorem 2 and conservativity of $L^{R}$ over $L(\backslash)$.

The other inclusion follows from Propositon 4.8: if we replace every type $C$ with $\operatorname{tr}(C)$ in an $\mathrm{L}^{\mathrm{R}}$-grammar, we obtain an L-grammar (since $\mathrm{L}^{\prime}$ and L differ only in the set of primitive types) generating the same language, and this language is now context-free by Theorem 3 .

Fragments of $\mathrm{L}^{\mathrm{R}}$ with restricted sets of connectives and/or primitive types are defined in the same way as for $L$.

Theorem 30. Derivability problems for $\mathrm{L}^{\mathrm{R}}\left(\backslash ; p_{1}\right)$, $\mathrm{L}^{\mathrm{R}}$, and all the calculi between them (fragments of $\mathrm{L}^{\mathrm{R}}$ and conservative extensions of $\mathrm{L}^{\mathrm{R}}\left(\backslash ; p_{1}\right)$ ) are NP-complete.

Proof. The derivability problem for $\mathrm{L}^{\mathrm{R}}$ is in the NP class due to Proposition 4.8 and the fact that the derivability problem for L is in NP.

NP-completeness of the derivability problem for $\mathrm{L}^{\mathrm{R}}\left(\backslash ; p_{1}\right)$ follows from the equivalence $B / A \leftrightarrow_{\mathrm{I}^{\mathrm{R}}}\left(A^{\mathrm{R}} \backslash B^{\mathrm{R}}\right)^{\mathrm{R}}$, that reduces derivability in $\mathrm{L}\left(\backslash, / ; p_{1}\right)$ to derivability in $\mathrm{L}^{\mathrm{R}}\left(\backslash ; p_{1}\right)$. The derivability problem for $\mathrm{L}\left(\backslash, / ; p_{1}\right)$ is NP-complete (Theorem 25).

## Bibliography

[1] A. V. Aho, R. Sethi, J. D. Ullman. Compilers: principles, techniques and tools. Reading, Mass.: Addison-Wesley, 1985.
[2] Y. Bar-Hillel, C. Gaifman, E. Shamir. On the categorial and phrase-structure grammars. Bulletin of the Research Council of Israel, Section F. 9F, 1960. - P. 1-16.
[3] W. Buszkowski. Compatibility of a categorial grammar with an associated category system. Zeitschrift für mathematische Logik und Grundlagen der Mathematik. 28, 1982. - P. 229-238.
[4] W. Buszkowski. Some decision problems in the theory of syntactic categories. Zeitschrift für mathematische Logik und Grundlagen der Mathematik. 28, 1982. P. 539-548.
[5] W. Buszkowski. The equivalence of unidirectional Lambek categorial grammars and context-free grammars. Zeitschrift für mathematische Logik und Grundlagen der Mathematik. 31, 1985. - P. 369-384.
[6] B. Carpenter. Type-logical semantics. - Cambridge, Mass.: The MIT Press, 1997.
[7] N. Chomsky. Three models for the description of language. IRE Transactions on Information Theory. I T-2, No. 3, 1956. - P. 113-124.
[8] M. Dekhtyar, A. Dikovsky. Generalized categorial dependency grammars. Trakhtenbrot/Festschrift, ed. by A. Avron et al. Lecture Notes in Computer Science, Vol. 4800. Berlin etc.: Springer, 2008. - P. 230-255.
[9] T. A. D. Fowler. A polynomial time algorithm for parsing with the bounded order Lambek grammars. The Mathematics of Language: proceedings of the 10-th and 11th Biennial Conference, MOL 2010 and MOL 2011, ed. by Ch. Ebert, G. Jäger and J. Michaelis. Lecture Notes in Computer Science, Vol. 6149. Berlin etc.: Springer, 2009. - P. 36-43.
[10] T. A. D. Fowler. Parsing CCGBank with the Lambek calculus. Parsing with Categorial Grammar Workshop, ESSLLI 2009. Bordeaux, 2009. - P. 38-42.
[11] J.-Y. Girard. Linear Logic. Theoretical Computer Science, 50, No. 1, 1987. - P. 1102.
[12] Ph. de Groote. A dynamic programming approach to categorial deduction. Proceedings of the 16th International Conference on Automated Deduction, ed. by H. Ganzinger. Lecture Notes in Artificial Intelligence, Vol. 1632. Berlin etc.: Springer, 1999. - P. 1-15.
[13] H. Hendriks. Studied flexibility. ILLC Dissertation Series, DS-1993-05, Amsterdam, 1993.
[14] G. Jäger. On the generative capacity of multi-modal categorial grammars. Research on Language and Computation, 1, No. 1-2, 2003. - P. 105-125.
[15] M. Kanazawa. The Lambek calculus enriched with additional connectives. Journal of Logic, Language and Information, 1, 1992. - P. 141-171.
[16] J. Lambek. The mathematics of sentence structure. American Mathematical Monthly. 65, No. 3, 1958. - P. 154-170.
[17] J. Lambek. Deductive systems and categories II: Standard constructions and closed categories. Category Theory, Homology Theory and Their Applications I, ed. by P. Hilton. Lecture Notes in Mathematics, Vol. 86. Berlin: Springer, 1969. - P. 76-122.
[18] J. Lambek. From categorial grammar to bilinear logic. Substructural Logics, ed. by K. Došen and P. Schroeder-Heister. Studies in Logic and Computation, Vol. 2. Oxford: Clarendon Press, 1993. - P. 128-139.
[19] F. Métayer. Polynomial equivalence among systems LLNC, $\operatorname{LLNC}_{a}$ and $\operatorname{LLNC}_{0}$. Theoretical Computer Science. 227, No. 1, 1999. - P. 221-229.
[20] M. Moortgat. Multimodal linguistic inference. Journal of Logic, Language and Information, 5, No. 3-4, 1996. - P. 349-385.
[21] M. Moortgat. Categorial type logic. Handbook of Logic and Language, ed. by J. van Benthem and A. ter Meulen. Elsevier, 1997.
[22] G. Morrill. Categorial grammar: logical syntax, semantics and processing. Oxford: Oxford University Press, 2011.
[23] M. Pentus. Equivalent types in Lambek calculus and linear logic. Preprint No. 2 of the Department of Mathematical Logic, Steklov Mathematical Institute, Series Logic and Computer Science, Moscow, 1992.
[24] M. Pentus. Free monoid completeness of the Lambek calculus allowing empty premises. Proceedings of Logic Colloquium '96, ed. by J. M. Larrazabal, D. Lascar and G. Mints. Lecture Notes in Logic, Vol. 12. Berlin etc.: Springer, 1998. - P. 171-209.
[25] M. Pentus. Lambek calculus is NP-complete. Theoretical Computer Science, 357, No. 1-3, 2006. - P. 186-201.
[26] M. Pentus. A polynomial-time algorithm for Lambek grammars of bounded order. Linguistic Analysis. 36, No. 1-4, 2010. - P. 441-471.
[27] M. Pentus. Complexity of the Lambek calculus and its fragments. Advances in Modal Logic, Vol. 8, ed. by L. Beklemishev, V. Goranko and V. Shehtman. - London: College Publications, 2010. - P. 310-329.
[28] D. N. Yetter. Quantales and noncommutative linear logic. Journal of Symbolic Logic, 55, No. 1, 1990. - P. 41-64.
[29] V. A. Minina. Completeness of the Lambek syntactic calculus with the involution operation (in Russian). Diploma paper, Dept. of Math. Logic and Theory of Algorithms, Moscow State University, 2001.
[30] M. R. Pentus. Lambek calculus and formal grammars (in Russian). Fundamental and Applied Mathematics, 1, No. 3, 1995. - P. 729-751.
[31] M. R. Pentus. Completeness of the Lambek syntactic calculus (in Russian). Fundamental and Applied Mathematics, 5, No. 1, 1999. - P. 193-219.
[32] Yu. V. Savateev. Algorithmic complexity of fragments of the Lambek calculus (in Russian). C. Sc. (Ph. D.) Thesis, Moscow State University, 2009.
[33] A. N. Safiullin. Derivability of admissible rules with simple premises in the Lambek calculus. Moscow University Math. Bulletin. 62, No. 4, 2007. - P. 168-171.

## Author's Papers

[34] S. L. Kuznetsov. Lambek calculus with one division and one primitive type permitting empty antecedents. Moscow University Math. Bulletin. 64, No. 2, 2009. P. 76-79.
[35] S. Kuznetsov. Lambek grammars with one division and one primitive type. Logic Journal of the IGPL. 20, No. 1, 2012. - P. 207-221.
[36] S. L. Kuznetsov. Lambek calculus with a unit and one division. Moscow University Math. Bulletin. 66, No. 4, 2011. - P. 173-175.
[37] S. L. Kuznetsov. Lambek calculus with the reversal operation (in Russian). VINITI preprint No. 152-V2012, 17.04.2012.
Revised English version: S. Kuznetsov. Lambek calculus with the reversal operation. Logical Aspects of Computational Linguistics, ed. by D. Béchet and A. Dikovsky. Lecture Notes in Computer Science 7351, 2012. - P. 151-160.

