Preface

FG-2006, the 11th conference on Formal Grammar, was held in Malaga, Spain in July 2006. This year's conference included 12 contributed papers covering, as usual, a wide range of topics in formal grammar. In addition to those papers, this volume includes also the abstracts of two invited talks by Josef van Genabith (Dublin City University) and Laura Kallmeyer (Universität Tübingen).

The twenty four submissions to the conference were reviewed by members of the Program Committee; we are grateful to all of them for their help in making the conference a success: Anne Abeille (Paris 7, FR), Pierre Boullier (INRIA, FR), Gosse Bouma (Groningen, NL), Chris Brew (Ohio State, US), Wojciech Buszkowski (Poznan, PL), Miriam Butt (Universität Konstanz, DE), Alexander Clark (Royal Holloway University, UK), Berthold Crysmann (DFKI, DE), Philippe de Groote (LORIA, FR), Denys Duchier (LORIA, FR), Tim Fernando (Trinity College, IE), Annie Foret (IRISA - IFSIC, FR), Nissim Francez (Technion, IL), Gerhard Jaeger (University of Bielefeld, DE), Aravind Joshi (UPenn, US), Makoto Kanazawa (National Institute of Informatics), Stephan Kepser (Tuebingen, DE), Alexandra Kinyon (University of Pennsylvania, US), Geert-Jan Kruijff (DFKI, DE), Shalom Lappin (King’s College, UK), Larry Moss (Indiana, US), Stefan Mueller (Universitäet Bremen, DE), Mark-Jan Nederhof (Max Planck Institute for Psycholinguistics, NL), James Rogers (Earlham College, US), Ed Stabler (UCLA, US), Hans Joerg Tiede (Illinois Wesleyan, US), Jesse Tseng (LORIA, FR), Willemijn Vermaat (Utrecht, NL), Anssi Yli-Jyrä (Helsinki, FI).

We are indebted to all the authors who submitted papers to the meeting, and to all participants in the Conference. On behalf of the Organizing Committee, which consisted of Paola Monachesi, Gerald Penn, Giorgio Satta and Shuly Wintner, I am happy to present this volume.

Shuly Wintner, February 2007
Contents

1 Constraint-based compositional semantics in lexicalized tree adjoining grammars 1
Laura Kallmeyer

2 Parsing and generation with treebank-based probabilistic LFG resources 5
Josef van Genabith

3 Treating clitics with minimalist grammars 9
Maxime Amblard

4 Linear grammars with labels 21
Houda Anoun & Alain Lecomte

5 P-TIME decidability of NL1 with assumptions 35
Maria Bulińska

6 Program transformations for optimization of parsing algorithms and other weighted logic programs 45
Jason Eisner and John Blatz

7 On theoretical and practical complexity of TAG parsers 87
Carlos Gómez-Rodríguez, Miguel A. Alonso, Manuel Vilares

8 Properties of binary transitive closure logics over trees 103
Stephan Kepser

9 Pregroups with modalities 119
Aleksandra Kislak-Malinowska
10 Simpler TAG semantics through synchronization
Rebecca Nesson and Stuart Shieber

11 Encoding second order string ACG with deterministic tree walking transducers
Sylvain Salvati

12 Sideways without copying
Edward P. Stabler

13 English prepositional passives in HPSG
Jesse Tseng

14 Linearization of affine abstract categorial grammars
Ryo Yoshinaka

15 List of contributors

Proceedings of FG-2006
Abstract

This talk presents a framework for LTAG semantics that computes semantics based on the LTAG derivation trees such that semantic computation consists of feature unifications parallel to those performed in Feature-Based TAG (FTAG). We show that this framework has sufficient expressive power to deal with a large range of seemingly problematic phenomena, namely quantifier scope, raising verbs, bridge verbs and nested quantificational NPs. Finally, a compositionality proof is sketched for this framework that relies on the fact that the derivation tree locally determines both, syntactic and semantic composition.\textsuperscript{1}

Keywords Lexicalized Tree Adjoining Grammars, computational semantics, compositionality, scope semantics, under-specification

1.1 Lexicalized Tree Adjoining Grammar (LTAG)

LTAG is a tree-rewriting formalism. An LTAG consists of a finite set of elementary trees associated with lexical items. From these trees, larger trees are derived by substitution (replacing a leaf with a new tree) and adjunction (replacing an internal node with a new tree). LTAG derivations are represented by derivation trees that record the way the elementary trees are put together. A derived tree is the result of carrying out the substitutions and adjunctions. Each edge in the derivation tree stands for an adjunction or a substitution. The elementary trees encapsulate all syntactic/semantic arguments of the

\textsuperscript{1}The work presented here can be found in Kallmeyer and Romero (2007) (for the framework and the scope analyses) and Richter and Kallmeyer (2007) (for the compositionality proof).
lexical anchor. They are minimal in the sense that only the arguments of the anchor are encapsulated, all recursion is factored out. Because of this, substitutions and adjunctions roughly correspond to combinations of a predicate with one of its arguments. Consequently, they determine semantic composition and therefore we compute LTAG semantics on the derivation tree.

1.2 LTAG Semantics with Semantic Unification

In our approach, each elementary tree is linked to a pair consisting of a semantic representation and a semantic feature structure description. These feature structure descriptions are used to compute assignments for variables in the representations using conjunction and additional equations introduced depending on the derivation tree.

The semantic representations consist of a set of labeled Ty2 formulas and a set of scope constraints of the form $x \geq y$ where $x$ and $y$ are propositional labels or propositional meta-variables. $x \geq y$ signifies that $y$ is a component of the term $x$. Meta-variables indicate that terms have not been specified yet. The assignment computed based on the feature structure descriptions specifies values for some of the meta-variables in the semantic representations while leaving some of them open. This allows for under-specified representations for scope ambiguities.

1.3 Scope Phenomena

In the talk we present analyses for the scope ambiguities exemplified in (1)–(5):

(1) Exactly one student admires every professor
   $\exists > \forall, \forall > \exists$

(2) John seems to have visited everybody
   $\text{seem} > \forall, \forall > \text{seem}$

(3) Three girls are likely to come
   $\text{three} > \text{likely}, \text{likely} > \text{three}$

(4) Mary thinks John likes everybody
   $\text{thinks} > \text{everybody}, *\text{everybody} > \text{thinks}$

(5) Two policemen spy on someone from every city
   $\forall > \exists > 2$ (among others), $*\forall > 2 > \exists$

Our analysis models the differences in scope behavior as follows:

1. Quantifiers scope within a kind of scope window delimited by an upper boundary maxs and a lower boundary mins, no matter where they attach inside a finite clause.
2. Operators on the verbal spine such as adverbs, raising verbs and bridge verbs take scope where they attach, i.e., among such operators, the attachment order specifies the scope order.

3. Adverbs and raising verbs are not concerned with the \textit{maxs–mins} scope window. Therefore, quantifiers can scopally interleave with them.

4. Bridge verbs embed a finite clause and in particular, they embed the \textit{maxs} limit of this clause. Therefore they block quantifier scope.

5. The maximal scope of a quantifier embedded in a quantificational NP is the proposition of the embedding quantifier. Therefore, if it scopes over the embedding quantifier, then this has to be immediate scope (no other quantifier can intervene).

1.4 Compositionality

At first sight, feature logic-based computational semantics systems such as LTAG do not seem compatible with a notion of compositionality. The derived trees clearly do not determine the meaning of a phrase in a compositional way. However, a crucial property of LTAG is that the derivation process (i.e., the process of syntactic combination) can be described by a context-free structure, namely the derivation tree. (This is why LTAG is mildly context-sensitive.) The way our LTAG semantics framework is defined, this context-free structure also specifies the process of semantic combination. In other words, we can define semantic denotations for the nodes in the derivation tree in such a way that the semantic denotation of a node depends only the denotations of the daughters, the semantic representation from the lexicon chosen for this node and the way the daughters combine with the mother. In this sense, LTAG semantics is compositional.

References


Treebank-based acquisition of “deep” grammar resources is motivated by the “knowledge acquisition bottleneck” familiar from other traditional, knowledge intensive, rule-based approaches in AI and NLP, following the “rationalist” research paradigm. Deep grammatical resources have usually been hand-crafted Butt et al. (2002), Baldwin et al. (2004). This is time consuming, expensive and difficult to scale to unrestricted text. Treebanks (parse-annotated corpora) have underpinned an alternative “empiricist” approach: wide-coverage, robust probabilistic grammatical resources are now routinely extracted (learned) from treebank resources Charniak (1996), Collins (1997), Charniak (2000). Initially, however, these resources have been “shallow”. More recently, a considerable amount of research has emerged on treebank-based acquisition of deep grammatical resources in the TAG, HPSG, CCG and LFG grammar formalisms. This talk provides an overview of research on rapid treebank-based acquisition of wide-coverage, robust, probabilistic, multilingual LFG resources. Grammar and lexicon acquisition O’Donovan et al. (2005) is based on an automatic LFG f-structure annotation algorithm Burke et al. (2004a), Burke (2006). I show how the acquired LFG resources can be used in wide-coverage, robust parsing Cahill et al. (2004) and generation Cahill and van Genabith (2006). I provide an overview of ongoing research on the induction of Chinese Burke et al. (2004b), Japanese, Arabic, Spanish, French and German Cahill et al. (2005) treebank-based LFG resources.
I briefly compare our LFG work with similar research on the treebank-based acquisition of HPSG Miyao et al. (2003) and CCG Hockenmaier and Steedman (2002) resources.

References


Treating clitics with minimalist grammars

Maxime Amblard

Abstract

We propose an extension of Stabler’s version of clitics treatment for a wider coverage of the French language. For this, we present the lexical entries needed in the lexicon. Then, we show the recognition of complex syntactic phenomena as (left and right) dislocation, clitic climbing over modal and extraction from determiner phrase. The aim of this presentation is the syntax-semantic interface for clitics analyses in which we will stress on clitic climbing over verb and raising verb.

Keywords Minimalist Grammars, syntax-semantic interface, λ-calculus, clitics.

Minimalist Grammars (MG) is a formalism which was introduced in Stabler (1997), based on the Minimalist Program, Chomsky (1995). The main idea which is kept from the Minimalist Program is the introduction of constituent move in the formal calculus. Such a “move” operation introduces flexibility in a system which seems to be like Categorial Grammars (CG). We try to recover the correspondence in CG, between syntactic structures and logical forms (interpretative level of the sentence).

This formalization introduces constraints on the use of move rules, and by this way makes the syntactic calculus decidable. These grammars are lexicalized and all steps of the analysis are triggered by the information extracted from the lexicon: from a sentence, it selects a subset of words. To each word corresponds a sequence of features, and it is the first element of the sequence in the derivation which triggers the next rules.

An advantage of this system is that the structure of the calculus is constant. The coverage of the grammar is extended by adding new elements to
the lexicon, never by adding new structural rules. The structural system of these grammars contains only two kinds of rules: move and merge (but extensions exist for both). We refer the reader to Stabler’s articles and others for presentation of the use of MG, Stabler (1997), Vermaat (1999).

Clitics are the normal form for pronoun in romance language. The syntactic and semantic behavior of clitics in these languages are complex. For French, clitics often climb over auxiliary verb. Ed Stabler proposes in Stabler (2001) a partial lexicon for French clitics recognition and analysis.

We propose here to extend this lexicon to several well-known linguistic problems. These problems interfere at different levels of analysis. Subject raising is typically a semantic question whereas the clitic climbing over modals is a syntactic question. We propose a new lexicon for its syntactic analysis and then we will show how our semantic interface solves semantic questions.

We use the description of clitics proposed by Perlmutter in Perlmutter (1971). He proposes a filter to recognize the right order of clitics for romance languages, from where we extract the sub-filter:

```
[je/tu/···]|ne|[me/te/se/···]|le/la/les/···]|lui/leur]|y/en].
[nominative | negative | reflexive | accusative | dative | locative | genitive].
```

In the first part, we propose an extension of Stabler’s version of clitics treatment for a wider coverage of the French language. For this, we will present the lexical entries needed in the lexicon. Then, we will show the recognition of complex syntactic phenomena as (left and right) dislocation, clitic climbing and extraction from determiner phrase. The aim of this presentation is the last part: the syntax-semantic interface for clitics analyses in which we will stress on clitic climbing over verb and raising verb.

### 3.1 Lexicon for French clitics

#### 3.1.1 Stabler analysis

Stabler’s works on clitics are inspired by Sportiche Sportiche (1992), who proposes the following treatment:

- Clitics are not elements moved from position XP*, but are co-referent with this position. The clitics appearing in the structure bear all the features their co-referring XP would bear. Furthermore, clitics do not form an autonomous syntactic object, but they are built into a unit with some host.

In this work, two parts in the cliticization are distinguished. The first one is an empty element which takes an argumental position from the verb. The second is the phonological treatment of the unit — the clitic in the surface structure.

We introduce lexical entries which are phonologically empty but carry special features which need to be unified with features of the phonological part
of the clitic. The two different parts are connected by a move operation. If just one of these items occurs in the sentence, derivation fails.

We sum up this treatment in the derivation as follows. The annotation recalls the main feature of the word and the annotation on the $e$ recall the word which 'exist trace':

\[(6) \text{ donne } e_F \rightarrow \text{Jean } _k \text{ la donne } e_lu \]
\[
\begin{align*}
\text{t} & \quad \text{Jean } _k \text{ la donne } e_lu \\
\text{Jean} & \quad \text{tJean } _k \text{ la donne } e_lu. \\
\text{John} & \quad \text{tJohn it gives } e_lu.
\end{align*}
\]

In more details, the derivation is the following:

**Derivation 1** Derivation of the simple French sentence: Jean la donne.

Lexicon:

<table>
<thead>
<tr>
<th>jean</th>
<th>D -k</th>
<th>donne</th>
<th>V</th>
<th>$\epsilon$</th>
<th>=T C</th>
<th>$\epsilon$</th>
<th>=&gt;V=D +k =D v</th>
<th>$\epsilon$</th>
<th>D -k -G</th>
</tr>
</thead>
</table>

Derivation step by step:

1. selection of lexical entry : \{ donne :: V\}
2. selection of lexical entry : \{ $\epsilon$ :: =>V =D +k =D v\} (which adds the syntactic component to the verb).
3. head movement. This is a merge between the two previous element where the phonological part of the argument moves to the phonological part of the head.
4. selection of lexical entry : \{ $\epsilon$ :: D -k -G\}. This is the empty argumental verb position.
5. merge.
6. There is a licensee “k” in first position, a move operation is triggered. After this step, the derivation tree is :

\[
\begin{align*}
\epsilon :: -G & \quad < \\
\epsilon :: =D v & \quad / \text{donne/}
\end{align*}
\]

7. selection of lexical entry : \{Jean :: D -k\}.
8. merge.
9. selection of lexical entry : \{la :: =v +G Acc3\}, the clitic takes part in the derivation.
10. merge.

11. move : the feature in the empty argument of the verb and the feature in the clitic are canceled.

12. selection of lexical entry : $[\epsilon ::=\text{Acc}3+KT]$ — to the end of the derivation.

13. merge.

14. move : resolution of nominative case :

```
   >
   |
  jean :: /Jean/  
   |
   |
   |
   |
   |
  /la/ ::  
   |
   |
  /donne/ ::
   |
   |
  \epsilon/ ::
```

15. selection of lexical entry : $[\epsilon ::=\text{T}C]$ — empty “complement” position.

16. merge ; end of the derivation with feature ‘c’ : acceptance.

In his presentation, Stabler proposes a lexicon for accusative, dative and reflexive clitics recognition. He ensures the right order with several verbal types. The analysis is driven by the head and the next clitics to introduce will have to be assigned verbal type as they occur in the Perlmutter filter’s order. Stabler uses the SMC — shortest move condition — to exclude the use of a reflexive and an accusative clitics in the same sentence.

### 3.1.2 Extension: genitive, oblique and nominative clitics

We can extend this first approach of French clitics treatment to other cases, in particular genitive, oblique and nominative. This section will present the lexical entries and the process of acceptance of derivations.

We call “state of a verb” the basic type of the head currently handled. For example, if a verb has a accusative clitic its type will be “Acc”.

For genitive and oblique clitics, we just add in the lexicon two new empty argumental positions and a list of possible types for each clitic.

In a first time, we introduce a new verbal type for beginning the cliticization and another where the cliticization is finished. We call them “clitic” and “endclitic”.

Treating clitics with MG / 13

Following the Perlmutter filter Perlmutter (1971), the first clitic we have to treat for keeping the right order is the genitive one. We add a genitive state which is connected to the “clitic” state. The verbal state passes to the genitive state by means of a lexical entry the phonological form of which is “en” and carries a licensee feature “en”:

\([en] :: [clitic \langle, +EN, genitive].\]

From this state we pass to all the other states of the cliticization, for example:

\([le] :: [genitive \langle, +G, acc].\]

and if there is only a genitive clitic, we use phonologically empty entry to pass to the end of the cliticization.

\([] :: [genitive \langle, endclitic].\]

The “oblique” clitics are treated in the same way, except that from “oblique” it is impossible to go back to “genitive”. All lexical entries of this type have a “y” phonological form.

\([y] :: [clitic \langle, +Y, oblique].\]
\([y] :: [genitive \langle, +Y, oblique].\]

In the same way, from oblique we can pass to other possible clitic states, as for example:

\([le] :: [oblique \langle, +G, acc].\]
\([leur] :: [oblique \langle, +F, dat].\]
\([] :: [oblique \langle, endclitic].\]

The nominative case is treated the same way. But the use of this procedure to add new clitic treatment is quadratic in the number of lexical entries. For the nominative pronoun, a discussion could be opened around its clitic state. We consider here that they are clitics.

Another discussion about negative form could rise around the status of the negation marker whose position is after the pronoun.

For the moment, we do not treat the negative form in a right way so we will not include it in this presentation, but we assume that the treatment of nominative clitics is outside the clitic cluster. All the phonological pronoun entries take a verbal form in “endclitic” state and give a new verbal form in “Nom” (inative) state.

We add an empty verb argument which must be included in the derivation before the clitic treatment:

\([] :: [d, −Subj, −case].\]

The sketch of the analysis is:

- la donne \(ε−Nom \)
We add in the lexicon a basic feature “Nom” and the lexical entries of the nominative pronouns, for example:

\[ \text{[je]} :: \{endclitic, +Subj, Nom\}. \]

\[ \text{[nous]} :: \{endclitic, +Subj, Nom\}. \]

The derivation continues with a phonologically empty entry at the end of the derivation.

\[ \text{[]} :: \{nom, +case, t\}. \]

### 3.2 Recognition of complex phenomena

This treatment of French clitics is simple and can be integrated easily into a larger analysis.

**climbing over modal**

We treat the clitic climbing over the whole verbal cluster in particular over modal.

The modal is combined with the verb in the inflection step. The inflection is treated with head movement and all clitics take their own place after this treatment.

If there are words which must be inserted between the verb and the modal — for sentences with adverbs — we first build the verbal constituent after which we treat the clitics. In this situation, the clitics could climb over the verb constituent or stand after.

For example, in French we can analyze a sentence as:

**(7)** Je l’ai vu.

I him have seen.
I have seen him.

by building the constituent ai vu. We can extend to sentences with inserted word: “Je l’ai souvent vu” / “I have often seen him” with a derivation as :

**(8)** ai souvent vu \(\epsilon\)→ \(\text{F} \)

\[
\begin{align*}
\text{Je} & \text{Nom} \rightarrow \text{Je} \text{l’ai souvent vu} \epsilon \rightarrow \text{Je} \text{l’ai souvent vu} \epsilon \rightarrow \text{Je} \text{l’ai souvent vu} \epsilon \\
\text{Je} & \text{Nom} \rightarrow \text{Je} \text{l’ai souvent vu} \epsilon \rightarrow \text{Je} \text{l’ai souvent vu} \epsilon \\
\text{Je} & \text{Nom} \rightarrow \text{Je} \text{l’ai souvent vu} \epsilon \rightarrow \text{Je} \text{l’ai souvent vu} \epsilon \\
\end{align*}
\]

I it often seen
I often saw him.
Clitic can be a direct recovery of a non-“empty verbal argument”, for example in case of nominal dislocation.

There is a non empty verbal argument which must be extracted from the main sentence and become an indirect argument of the verb.

We build a verb with an “argument which must be extracted” — a determiner phrase (DP) — must be outside the main sentence. This state is introduced by a pause or comma. It modifies the determiner phrase in two different ways which depend on the side of the extraction:

- it adds a licensee for the left dislocation and cliticization.
- it adds a licensee for cliticization (and nothing for right dislocation).

The main problem is to include in the sentence the right part which will be replaced by the clitic.

Left dislocation: the DP is extracted from the sentence, placed in first position and recovered by a clitic.

\[
(9) \text{Marie le voir trop ce type, } \rightarrow \text{Ce type, Marie le voir trop.}
\]

That guy, Marie him sees too much.

Lexical entry of modifier of DP.

\[
[. .] :: [= \text{-}d, d, -H, \text{disloc}].
\]

Remark that we use comma to caring this treatment, but it can be through an empty lexical entry. The analysis would be the same. The comma will be placed after the DP by a head movement. The first licensee will be canceled with the licensor of the clitic and the second with another entry that we must add in the classical “comp” entry (this last entry is used to finish the derivation).

\[
[. .] :: [t, c, +DIS LOC].
\]

Right dislocation: In this case the determiner phrase is placed at the end of the sentence. For the homogeneity of the mechanism, we add a licensee of recovered by a clitic, and another for the extraction at the end of the sentence.

\[
[. .] :: [d <-, d, -H, \text{disloc}].
\]

The “comp” phase uses a weak move which lets the phonological form of the constituent in its place — here, at the end of the sentence.

\[
(10) \text{Marie le voir trop, ce type, } \rightarrow \text{Marie le voir trop, ce type.}
\]

Marie him sees too much, that guy.

This extraction seems to be very similar to questions: in questions, an argument of the verb is extracted to take another position in the surface level of
the sentence.

**Extraction from DP**

With the same kind of mechanism, we can extract an argument of any constituent. The determiner phrase can be complex and we extract an argument of the DP. For example:

   Peter of-it sees the end — Peter sees the end of the movie.

We build “la fin ε-\text{en}” and the cliticization allowed the extraction of the genitive. “Pierre en voit la fin.”

**Raising verb**

Raising verbs are verbs where one of the arguments is a verb and one of the other arguments is shared by both verbs, like in the sentence:

(12) Il semble le lui donner.
   He seems it him give.
   He seems give it to him.

where the pronoun “Il” is subject of the two verbs “semble” and “donner”.

The second verb must be in infinitive form.

In this case, the sentence has the following structures:

[subject raising\verb clitic infinitive\verb ].

A raising verb takes as an argument a verb in infinitive form, with a special inflection “infinitive”, and without subject. The infinitive inflection has the lexical entry:

[-inf]:[=>v, verbe].

“verbe” is the feature needed before starting the clitic treatment. A verbal form gets a “verbe” type after the verb receives its inflection.

The raising verb selects such a “verb”, then a DP subject and then becomes a VP of type “raisingv” which means a VP which has not yet received the inflection feature and will be able to receive new clitics (in particular pronoun).

For example:

[semble]:[=verbe, =d, raisingv].

This verb should receive its inflection and its subject. It follows this mechanism until the end of the derivation:

* semble la répare-inf
* semble -ε la répare-inf
* Je semble -ε la répare-inf
* I seem -ε it repara-inf
I seem repair it

### 3.3 Semantic interface

#### 3.3.1 How to use the syntax/semantic interface

From a sentence, we build a formula of higher order logic which represents its propositional structure. We associate to each lexical entry a $\lambda$-term and to each syntactic rule an equivalent semantic rule. We assume that the syntactic analysis drives the semantic calculus.

$\lambda$-terms application occurs only when an element has no features. We assume the following functions:

\[
\text{feat}(x) = \begin{cases} 
1 & \text{if the number of feature of } x = 0 \\
0 & \text{else}
\end{cases}
\]

\[
\text{sem}(x, y) = \begin{cases} 
1 & \text{if } \text{feat}(x) = 1 \text{ or } \text{feat}(y) = 1 \\
0 & \text{else}
\end{cases}
\]

Syntactic and semantic synchronization: after any operation in the syntactic calculus, the semantic counterpart computes the $\text{sem}$ function and if $\text{sem}(x, y) = 1$, we perform the functional application of the two $\lambda$-terms. To known which application to perform, we look at the type of the semantic terms.

A semantic tree represents the semantic counterpart of the sentence. It is a tree where the leaves are the semantic part of the lexical entries and the inner nodes contain the $\lambda$-term built and the direction of the head (of the syntactic part). We use the following notation:

- breaker between direction head and $\lambda$-term: $\vdash$.
- application: @

Applications are carried out when syntax allows it, therefore when the function $\text{sem} = 1$ for one of the two terms. The following applications are possible:

\[
\begin{align*}
\text{if } \text{sem}(\lambda\text{-term 1}, \lambda\text{-term 2}) = 1 \quad &\text{else} \\
\vdash \lambda\text{-term 1} @ \lambda\text{-term 2} \\
\vdash \lambda\text{-term 1} \land \lambda\text{-term 2} \\
\vdash \lambda\text{-term 1} \land \lambda\text{-term 2}
\end{align*}
\]

If a move operation canceled the last feature, we represent it by a unary branch in the tree.

Remark. There are two different possibilities for the semantic calculus: either waiting for elements completely discharged either immediately perform the application. But both fail in different cases: immediate application fails in case of “late adjunction” and the other possibility fails in questions treatment. The right solution seems to be intermediate: it consists in determining
a subset of features which must be consumed before applications will be performed. For the moment, we choose the first possibility. Later on, we shall do differently but this only involve changes in the feat function.

### 3.3.2 Example of semantic treatment

#### Clitic semantics

We present a syntactic treatment of clitics in two different parts. One is phonologically empty and is the non empty argument of the verb, the other is syntactically empty but it is a phonological recovery of the first one. The semantical part of the clitic is in the argumental position and this is a free variable which must be bound in the context. The phonological recovery is an identity.

<table>
<thead>
<tr>
<th>lexical entries</th>
<th>syntactic form</th>
<th>semantic form</th>
</tr>
</thead>
<tbody>
<tr>
<td>la</td>
<td>dat &lt;= +G acc</td>
<td>Id</td>
</tr>
<tr>
<td>t(la)</td>
<td>p - case - G</td>
<td>x*</td>
</tr>
</tbody>
</table>

* Free variable, bound in the context — we could use the Bonato algorithm to determine how this variables are bounded Bonato (2006).

We briefly present a semantic tree for a clitic treatment:

(13) Jean la répare.

John it repairs.
John repairs it.

In the semantic tree of the part of the cliticization above, we do not represent the identity operator (except for the clitic one).

\[
\begin{align*}
\langle t(la) \rangle @ & \text{Infl} @ \text{donne, je} \\
\langle t(la) \rangle @ \text{Infl} & @ \text{donne, je, t(la)} \\
la :: \text{Id} & @ \langle t(la) \rangle \\
\text{Infl} & \langle \text{donne, je, t(la)} \rangle \\
\text{Infl} & \langle \text{donne, t(la)} \rangle & \text{je} \\
\text{Infl} & \langle \text{donne, t(la)} \rangle & \text{je} \\
\end{align*}
\]

The last part of the tree is built by a move which creates a link between the phonological part of the clitic and the argumental part.

#### Over raising verbs

For the semantic calculus, raising verbs are predicates which take a subject and an action as argument. They apply a variable at this action.

We present the analysis of the sentence:
(14) Je semble la réparer.
    I seem it repair.
    I seem repair it.

The λ-terms, semantic counter-part of lexical entries are:

<table>
<thead>
<tr>
<th>sembler</th>
<th>ΛSΛv.(seem v, S(v))</th>
</tr>
</thead>
<tbody>
<tr>
<td>Je</td>
<td>I</td>
</tr>
<tr>
<td>ελy</td>
<td>Y*</td>
</tr>
<tr>
<td>réparer</td>
<td>λ x λ x y . repair (y, x)</td>
</tr>
</tbody>
</table>

* this variable is bound in the context

The semantic counter part of the pronoun is a constant referring to the speaker “I”. The clitic subject climbs over the raising verb. It can be the subject of both verbs in the sentence due to the semantic structure of the raising verb. If the main verb of the sentence has a subject, the application will not introduce a new variable in the formula, else the main verb needs a variable which stands at the subject place. The raising verb involves this variable by duplication of its subject.

The syntactic analysis builds the following structure:

(I@(inflection@(seem@(la(infinite@repare)))))

which allows the computation of the formula: “la reparer”

λx.repair(x, Y)

and this term is applied to the raising verb: ΛSΛv.(seem v, S(v))

λv.seem(v, repair(v, Y))

At the end of the calculus, we construct the formula:

pres(seem(I, repair(I, Y)))

where Y is bound in the context.

This is the formula we want to construct for representing the propositional semantics of the sentence. The subject clitic syntactically climbs over the main verb, and semantically climbs over the two verbs.

### 3.4 Conclusion and future work

In this paper, we presented an extension of Ed Stabler’s propositions on French clitics in minimalist grammars. The new lexicon makes it possible to treat several other syntactic phenomena, the same way as clitic climbing, e.g. extraction from NP or right and left dislocation.

Then, we proposed a syntax-semantic interface for Minimalist Grammars. The aim of this calculus is to build a formula of higher order logic. The semantic calculus, λ-calculus, is driven by the syntactic one. We emphasize on
the way to recognize clitics and semantic implication of climbing with raising verbs.

For future work, we want to integrate the negation into the grammar. We consider that the neg-marker “ne” is a clitic and must be incorporated in the treatment of French clitics. There is another complex phenomenon to consider concerning with clitics in the imperative mode (and negation).

Other cases of raising verbs exist which are more complex, allowing several syntactic clitic climblings as in:

(15) Je la laisse le lui donner.
    I her let it (to) him give.
    I let her give it to him.

where clitics take place in different orders.

Moreover, we want to continue to model the semantic effect of clitics in sentences, in particular for interaction between quantifier scope and clitics, which can introduce ambiguities in sentences like:

(16) Je la laisse tous les lui donner.
    I her let all them him give.
    I let her gives all to him.

Acknowledgment

The writer like to thank Christian Retoré and Alain Lecomte for crucial supports and one of the anonymous TALN 2006 referees for important comments and examples reuse in this paper.

References


Linear grammars with labels

HOUDA ANOUN & ALAIN LECOMTE

Abstract
The purpose of this paper is to show that we can work in the spirit of Minimalist Grammars by means of an undirected deductive system called LGL enhanced with constraints on the use of assumptions. Lexical entries can be linked to sequences of controlled hypotheses which represent intermediary sites. These assumptions must be introduced in the derivation and then discharged in tandem by their proper entry which will therefore manage to find its final position: this allows to logically simulate move operation. Relevance of this formalism will be stressed by showing its ability to analyze difficult linguistic phenomena in a neat fashion.

Keywords Logicoal grammars, Minimalist Program, syntax/semantics interface, non-linear phenomena

4.1 Introduction
Type Logical Grammars (Lambek (1958), Moortgat (1997)) and Minimalist Grammars (Chomsky (1995), Stabler (1997)) are two thriving theories dedicated to natural language analysis. Each one has its intrinsic assets. In fact, the first framework is computationally attractive as it works compositionally and gives the semantics for free. While the second one is based upon a reduced number of rules guaranteeing processing efficiency (Harkema (2000)).

Despite their apparent differences, these theories share the same philosophy: they are both lexicalized and present universal sets of rules that allow to explain various linguistic phenomena in multitude of natural languages.

Our goal is to bridge the gap between Categorial and Minimalist Grammars by proposing a new logical formalism LGL (i.e. Linear Grammars with Labels) which captures Minimalist operations (i.e. merge and move) in a deductive setting. This match between logical framework and Minimalist Program proves to be fruitful as it gives a better understanding of the different
mechanisms involved in Minimalist derivations.

Lecomte, A. and Retoré, C. have already proposed a logical system that simulates Minimalist Grammars: Lecomte and Retore (2001). This latter system is built upon elimination rules for both the slashes and the tensor. The absence of any form of introduction rules leads to an efficient system. However, this restriction is not beneficial insofar as it violates the correspondence between syntactic types and semantic representations. In our new proposal, we want to keep a transparent interface between syntax and semantics by reintroducing abstraction rules which are applied in a controlled fashion.

Like Abstract Grammars and Lambda-Grammars (de Groote (2001) and Muskens (2003)), LGL grammars are based upon an undirected logical system which has two interfaces (syntactic-phonetic, syntactic-semantics) owing to Curry-Howard correspondence. A syntactic derivation is then a deductive proof of a given sequent built using appropriate inference rules. Both phonetic form and semantic representation result from \( \lambda \)-terms combination which is carried out in parallel with the syntactic derivation, since each deductive rule encapsulates a computational step within the simply typed \( \lambda \)-calculus.

The originality of LGL stems from the refinement introduced in hypothetical reasoning. Our model aims at preserving the advantages of this technique (e.g. dealing with unbounded dependencies) while constraining its use in order to reduce the size of the search space. Thus, instead of considering freely accessible logical axioms, our system is equipped with finite sequences of consumable controlled hypotheses which are attached to certain lexical entries that are expected to move. Such linked hypotheses represent original sites occupied by their associated entry in the D-structure (i.e. before the displacement operation). They should be introduced during the derivation and then abstracted at the same time by their proper entry which will consequently reach its target. In the case of overt constituent movement, intermediary positions occupied by non-pronounced variables will be systematically replaced by phonetically-empty traces.

In this paper, we will prove that move is a metaphoric notion which can be rigorously formalized using Logic. Moreover, we will show how to capture complex linguistic phenomena (e.g. binding, discontinuity) within LGL thanks to the combination between Logic power and Minimalist Program ideas.

4.2 Bases of LGL

4.2.1 Types & Terms

In this section, we survey the relevant bases inherent to LGL.

Following earlier proposal by Curry, HB. in Curry (1961) and other more recent research work: de Groote (2001), Muskens (2003), our system dis-
tistinguish between two fundamental levels of grammar. The first level is an abstract language (tectogrammar) which encapsulates universal principles. The second level is a concrete one which may contain a range of components (e.g. phenogrammar, semantics) used to encode cross-linguistic variation (e.g. word order, lexical semantics).

Our core logic operates on abstract syntactic types which are inductively defined as follows:

\[ \mathcal{T}(\mathcal{A}) := \mathcal{A} \mid \mathcal{T} \rightarrow \mathcal{T} \mid !\mathcal{T} \]

\(\mathcal{A}\) is a finite set of atomic types that contains usual primitives of minimalist grammars (e.g. \(e\) (sentence), \(d_{\text{acc}}\) (noun phrase with accusative case), \(d_{\text{nom}}\) (noun phrase with nominative case). Composite types are built using the linear implication \(\rightarrow\) and the exponential operator \(!\) introduced in Girard (1987).

Our framework supports a two-dimensional concrete level dealing respectively with phonetics and semantics. Therefore, we consider two kinds of concrete types, namely \(\Phi\)-types \((\mathcal{T}_\Phi)\) and \(\lambda\)-types \((\mathcal{T}_\lambda)\) whose definitions are the following:

\[ \mathcal{T}_\Phi := s \mid \mathcal{T}_\Phi \rightarrow \mathcal{T}_\Phi \]

\[ \mathcal{T}_\lambda := e \mid t \mid \mathcal{T}_\lambda \rightarrow \mathcal{T}_\lambda \]

The set \(\mathcal{T}_\Phi\) is composed of only one atomic type \(s\) which represents phonetic structures (structured trees), whereas \(\mathcal{T}_\lambda\) contains two primitives \(e\) (individuals) and \(t\) (truth values). Notice that composite \(\Phi\)-types are built upon linear implication \(\rightarrow\), whereas composite \(\lambda\)-types use intuitionistic implication \(\to\).

Both phonetic and semantic representation of expressions are easily defined owing to \(\lambda\)-calculus, thus leading to two sets of terms, namely \(\Phi\)-terms \(\Lambda_\Phi\) and \(\lambda\)-terms \(\Lambda_\lambda\). Let \(\Sigma\) be a finite set of phonetic constants and \(C\) a finite set of semantic constants. Let \(\mathcal{V}_\Phi\) (resp. \(\mathcal{V}_\lambda\)) be an infinite countable set of typed phonetic (resp. semantic) variables. The set \(\Lambda_\Phi(\Sigma)\) of well-typed linear \(\Phi\)-terms is inductively defined as follows:

1. \(\epsilon \in \Lambda_\Phi(\Sigma)\) and \(\epsilon\) is of type \(s^1\)
2. if \(\phi \in \Sigma\) then \(\phi \in \Lambda_\Phi(\Sigma)\) and \(\phi\) is of type \(s\)
3. if \((x_\Phi : t_\Phi) \in \mathcal{V}_\Phi\) then \(x_\Phi \in \Lambda_\Phi(\Sigma)\)
4. if \(s_1\) and \(s_2\) are \(\Phi\)-terms of type \(s\) then \(s_1 \cdot s_2 \in \Lambda_\Phi(\Sigma)\) and it is of type \(s\) (\(\cdot\) operator is used to combine phonetic structures, it is neither associative nor commutative)
5. if \(\phi_1\) and \(\phi_2\) are \(\Phi\)-terms of types \(t_1\) and \(t_1 \to t_2\) with no common free variable then \((\phi_1 \phi_2) \in \Lambda_\Phi(\Sigma)\) and is of type \(t_2\)
6. if \(x_\Phi\) is a variable of type \(t_1\), \(\phi_1\) a \(\Phi\)-term of type \(t_2\) and \(x_\Phi\) occurs free exactly once in \(\phi_1\) then \((\lambda x_\Phi \phi_1) \in \Lambda_\Phi(\Sigma)\) and has type \(t_1 \to t_2\)

\(^1\epsilon\) represents a phonetically empty element used for traces
\[ \Lambda_\Phi(\Sigma) \text{ is provided with the usual relation of } \beta\text{-reduction } \Rightarrow \text{ enhanced with two additional rewriting rules: } \phi_1 \cdot e \Rightarrow \phi_1 \text{ and } e \cdot \phi_1 \Rightarrow \phi_1. \]

On the other hand, the set \( \Lambda_\Lambda(C) \) of \( \Lambda \)-terms is defined using a simply typed \( \Lambda \)-calculus with two basic operations, namely intuitionistic application and abstraction.

Finally, let \( \tau_{lat} \) be a function which assigns a \( \lambda \)-type to each atomic abstract type (we assume for instance that: \( \tau_{lat}(c) = t \), \( \tau_{lat}(n) = e \rightarrow t \), \( \tau_{lat}(d \text{ case}) = e \)).

Two homomorphisms \( \tau_\Phi \) and \( \tau_\lambda \) are defined to link abstract types to concrete types as follows:

\[
\begin{align*}
\forall t \in A, \tau_\Phi(t) &= s \\
\tau_\lambda(t_1 \cdot t_2) &= \tau_\lambda(t_1) \Rightarrow \tau_\lambda(t_2) \\
\tau_\lambda(t_1) &= \tau_\lambda(t_2)
\end{align*}
\]

4.2.2 Lexical Entries & Controlled Hypotheses

We now introduce the notion of 2-dimensional signs which are the basic units managed by our system. Such signs are of the following form \((l_\Phi, l_\lambda) : ty\), where:

- \( ty \in T(A) \) (abstract type)
- \( l_\Phi \in \Lambda_\Phi(\Sigma) \) and \( l_\Phi \) is of concrete type \( \tau_\Phi(ty) \)
- \( l_\lambda \in \Lambda_\Lambda(C) \) and \( l_\lambda \) is of concrete type \( \tau_\lambda(ty) \)

We distinguish between three classes of signs, namely variable signs (when \( l_\Phi \in V_\Phi \) and \( l_\lambda \in V_\lambda \)), constant signs (when \( l_\Phi \in \Sigma \) and \( l_\lambda \in C \)) and compound signs (when \( l_\Phi \) or \( l_\lambda \) is a compound term).

These signs are used to define lexical entries. Lexical entries of \( LGL \) are proper axioms which can be coupled with prespecified sequences of controlled hypotheses. Such hypotheses will occupy intermediary sites, they should be introduced in the appropriate order and then discharged at the same time by their associated entry.

Lexical entries obey the syntax below:

\[ \vdash (a_\Phi, a_\lambda) : ty \Rightarrow l_{\text{hyps}} \]

where:

- \((a_\Phi, a_\lambda) : ty \) is a 2-dimensional sign.
- \( l_{\text{hyps}} = ([H_1 : t \Rightarrow H'_1 : t], ..., [H_k : t \Rightarrow H'_k : t]) \) is a sequence of controlled axioms of length \( |l_{\text{hyps}}| = k \), \( (\forall i \in [1.k], H_i \equiv (h_{\Phi_i}, h_{\Lambda_i}) \text{ and } H'_i \equiv (h'_{\Phi_i}, h'_{\Lambda_i})) \)

Lexical entries are classified in two groups: linked entries (when \( k > 0 \)) and free ones (when \( k = 0 \)). Linked entries are coupled with non-empty sequences of controlled hypotheses. Each hypothesis is encapsulated inside an axiom
Linear grammars with labels

\'(h_{st}, h_{st}) : t \vdash (h'_{st}, h_{st}) : t\' which can be either logical (if \( h_{st} = h'_{st} \)) or extra-logical (if \( h_{st} \neq h'_{st} \)). Extra-logical axioms are extremely useful since they represent pronounced variables or phonetically non-empty traces stemming from displacement (e.g. pronouns: he, her ...).

The abstract type \( ty \) of the lexical entry should verify the following specification:

1. if \( k = 0 \) then \( ty \) is an arbitrary abstract type
2. if \( k = 1 \) then \( ty = \tilde{t}_1 \rightarrow^{...} \rightarrow \tilde{t}_n \rightarrow (t' \rightarrow \tilde{t}' \rightarrow \tilde{t}'') \)
3. otherwise \( ty = \tilde{t}_1 \rightarrow^{...} \rightarrow \tilde{t}_n \rightarrow (! t' \rightarrow \tilde{t}' \rightarrow \tilde{t}'') \)

Intuitively, the second (resp. third) point above means that our lexical entry represents a constituent that needs to merge with exactly \( n \) (\( n \geq 0 \)) expressions of types \( t_1 \ldots t_n \) respectively, and then move once (resp. an unspecified number of times, e.g. cyclic move) to reach its final position.

Finally, a lexicon is nothing else but a finite set of lexical entries \( \{e_1, \ldots, e_n\} \).

Let us illustrate the previous definitions in a concrete example. If we assume that \( (\text{whom} \in \Sigma) \) and \( (\wedge \in C) \) then the phonetic behavior and the semantic representation of the relative pronoun \( \text{‘whom} \) can be modeled using the linked entry below:

\[ \vdash \left( \lambda \phi \lambda m. m \bullet (\text{whom} \bullet \phi(x)) \right) \colon (d_{acc} \rightarrow c) \rightarrow n \rightarrow n \rightarrow [X : d_{acc} \vdash X : d_{acc}] \]

Our entry is linked to one hypothesis which will occupy the initial position of \( \text{‘whom} \), namely the object of its relative clause (e.g. \( \text{(book) whom Noam wrote} \)). This assumption will be discharged afterwards by its related entry, thus guaranteeing the combination between the relative pronoun and its subordinate clause. Formal rules that manage this overt displacement will be set forth in the next section.

### 4.3 Logical simulation of Minimalism

#### 4.3.1 Inference rules

Let \( lex = \{e_1, e_2, \ldots, e_n\} \) be a lexicon. \( LGL \) grammar with lexicon \( lex \) is based upon a deductive logical system which deals simultaneously with two interfaces (syntactic-phonetic, syntactic-semantic).

Judgments of our calculus are sequents of the following form:

\[ \Gamma \vdash (l_\Phi, l_\lambda) : ty \, ; \, E \]

where:

- \( \Gamma \) the context is a finite multiset of 2-dimensional variable signs
- \( (l_\Phi, l_\lambda) : ty \) is a 2-dimensional sign
• $E$ is a finite multiset containing identifiers of all linked lexical entries that were used in the course of the derivation and whose associated assumptions are not yet discharged.

Variable signs included in the context $\Gamma$ correspond to controlled hypotheses that were introduced in the course of the derivation. Each hypothesis will be marked using a superscript ‘$\uparrow$’ which points at the lexical entry to which the assumption is attached (e.g. $x_i^\uparrow\Phi$: hypothesis linked to $e_i$ entry).

The first group of $\mathcal{LGL}$ inference rules are axioms which coincide with derivations’ leaves. Figure 1 shows axioms that our system supports:

\[
\begin{align*}
\frac{\vdash a_\Phi : ty; \text{ if } l = (\text{then } \emptyset \text{ else } \{e_i\})}{e_i = (\vdash a_\Phi : ty; l)} \quad \text{Lex} \\
\frac{\vdash a_\Phi : ty; l_{hyp}[j] = (x_\phi : A \vdash y_\phi : A) \quad x_i^\uparrow : A \vdash y_\phi : A; \emptyset}{e_i = (\vdash \_ \vdash l_{hyp}) \quad \text{Ctrl}}
\end{align*}
\]

**FIGURE 1** Axioms of $\mathcal{LGL}$($\text{lex}$)

Our core logic includes extra-logical axioms which are extracted from lexical entries owing to rule $\text{Lex}$. If the involved entry is linked, then its identifier is added to the multiset $E$. On the other hand, our system excludes the freely accessible identity axiom. Available axioms stem from controlled hypotheses which are coupled with linked lexical entries. These axioms can be introduced in the derivations by means of $\text{Ctrl}$ rule.

Linked entries in $\mathcal{LGL}$ can be attached to more than one controlled hypothesis. This specification has a very strong linguistic motivation. In fact, it can happen that a constituent occupies more than one intermediary site before reaching its target. Such phenomenon is illustrated for instance in the interrogative sentence ‘Which book did John file without reading it?’. In that case, the wh-element ‘which book’ occupied two positions before displacement (in the D-structure), namely the complement of the verb file and that of the infinitive without reading. After movement, the first position becomes empty while the second is occupied by a pronounced variable ‘it’. At the semantic level, both these sites of origin represent the same object.

To account for such non-linear phenomena within $\mathcal{LGL}$, we use the exponential $!$ whose behavior is described by the usual rules of linear logic (Girard (1987)). Figure 2 presents the derived rules which are relevant to our study.

The generic process that handles the management of controlled hypotheses can be summarized as follows. On the first hand, each assumption of type

---

2 For the sake of readability, we focus on the syntactic-phonetic interface
\[
\Delta, x_\phi^y : B \vdash y_\phi : A; \quad E_1 \quad \vdash L
\]

\[
\Delta, x_\phi^y : !B, y_\phi : !A; \quad E_1 \quad \vdash !L
\]

\[
\Delta, b_\phi^x : !B, y_\phi : !B \vdash u_\phi : A; \quad E_1 \quad \vdash !L
\]

**Figure 2** Relevant derived rules for \( ! \)

\( ty \) will get the decorated type \( !ty \) if it is related to a linked entry \( e_i \) which is attached to more than one controlled hypothesis. This transformation is carried out by means of \( !L \) rule. Intuitively, this means that a hypothesis which represents only one controlled assumption (i.e. of type \( ty \)) is a particular case of hypotheses that encapsulate at least one controlled assumption (i.e. of type \( !ty \)). On the second hand, contraction rule \( !L \) is applied to gather all the hypotheses linked to a specific entry \( e_i \) in one assumption. This will make it possible to abstract these hypotheses in tandem.

Now, the ground is well prepared to present our logical simulation of Minimalism. It is not difficult to simulate merge operation of Minimalist Grammars in a logical setting. In our case, it is nothing else but the direct \( \to \) elimination (\( \to E \), cf. Fig. 3) which merges two \( \Phi \)-terms (resp. \( \lambda \)-terms) by means of application operation.

\[
\Gamma \vdash f_\phi : A \to B; \quad E_1 \quad \Delta \vdash a_\phi : A; \quad E'_1 \quad \Gamma, \Delta \vdash (f_\phi a_\phi) : B; \quad E_1 \cup E'_1 \quad \to E
\]

\[
\Gamma \vdash f_\phi : (C \to D) \to B; \quad [e_i] \cup E_1 \quad \Delta, c_\phi^x : C \vdash d_\phi : D; \quad E'_1 \quad \Gamma, \Delta \vdash (f_\phi (\lambda c_\phi. d_\phi)) : B; \quad E_1 \cup E'_1 \quad \to IE \triangleleft
\]

**Figure 3** Behavior of \( \to \) connective

Move operation is logically captured thanks to the refined elimination rule \( \to IE \). This rule allows a constituent to reach its final position by simultaneously discharging its controlled hypotheses which occupied intermediary positions. Our logical formalization of move operation shares some ideas with Vermaat’s one in Vermaat (1999). In fact, we both consider this operation as the combination of two phases, namely a merge step and a hypothetical reasoning\(^3\) step (abstraction over sites of origin). Thus, the elements which are expected to move are assigned a higher order type \( (C \to D) \to B \)\(^4\). Such elements wait to merge with a constituent of type \( C \to D \), which results from the abstraction of the intermediary positions in the initial structure (of type \( D \)).

However, Vermaat proposal is encoded in a directional calculus: move operation is then captured using additional postulates which reintroduce struc-

\(^3\) The introduction rule of \( \to \) is not freely available, it is rather encapsulated inside \( \to IE \) rule

\(^4\) Vermaat considers only the case where \( D = B \)
atural flexibility in a controlled fashion. Our proposal is simpler as it is based
upon a flexible undirected calculus. Moreover, it makes it possible to limit
the operation of hypothetical reasoning used in displacement which is con-
strained to a specific amount of hypotheses explicitly given by the lexicon.

Rule ¬IE cannot be applied unless the pre-condition \( \hat{z} \) is verified: all
linked axioms coupled with the lexical entry \( e_i \) must be introduced in an ap-
propriate order (from the right to the left of \( l_{hyp} \) sequence) during the deriva-
tion of \( (\Delta, c'_\phi : C \vdash d_\phi : D; E'_1) \). Once these assumptions are abstracted, entry
\( e_i \) regains its final position and is automatically withdrawn from the multiset
of unstable lexical entries involved in the derivation.

To formalize the pre-condition \( \hat{z} \), we assume that each assumption \( x^V \) of
the context encapsulates a kind of history used to record some relevant data.
This additional parameter does not have any impact on our logical system.
It only ensures the efficiency of parsing by making the constraint \( \hat{z} \) easier
to check. The notation \( x^V[\sigma] \) is used when the history \( \sigma \) of the assumption
\( x^V \) is explicitly given. Otherwise, a function \( hist() \) can be applied to a given
hypothesis \( x^V \) to get its masked history.

Owing to the contraction rule \(!Lc\), each hypothesis \( x^V \) gathers a sub-set
of controlled hypotheses related to entry \( e_i \). The history of an assumption \( x^V \)
can then be encoded as a set of pairs of natural numbers. The first number
of each pair represents the index of an involved controlled hypothesis taken
from \( l_{hyp} \) sequence, while the second one is nothing else but the depth\(^5\) of
this hypothesis in the current bottom-up derivation.

Each deduction step updates the history of all assumptions included in the
context. For instance, \( Ctrl \) rule enables the introduction of a specific con-
trolled hypothesis of index \( j \) and initiates its history with the single pair \((j, 0)\).
On the other hand, rules of Fig.2 and Fig.3 increment\(^6\) the depth of the pre-
viously introduced controlled hypotheses. We show below two logical rules
enhanced with their explicit management of histories:

\[
\frac{e_i = (\vdash \neg \exists l_{hyp}) \quad l_{hyp}[j] = (x_0 : A \vdash y_0 : A)}{x^V[(j, 0)] : A \vdash y_0 : A; \emptyset \quad Ctrl}
\]

\[
\frac{\Delta, x^V[\sigma_1] : !B, y^V[\sigma_2] : !B \vdash u_\phi : A; E_1}{\Delta, b^V[\sigma_1^+ \cup \sigma_2^+] : !B \vdash u_\phi[x_\phi := b_\phi, y_\phi := b_\phi] : A; E_1 \quad \!Lc}
\]

\(^5\)The number of deduction steps between the introduction of the hypothesis and the current
state of the derivation

\(^6\)Incrementing operation is denoted by \((\ldots k_i d_i \ldots)^+ = (\ldots (k_i d_i + 1) \ldots)\)
Therefore, the side condition $\ddagger$ can be stated formally as follows:

$$\ddagger \iff \forall k, 1 \leq k \leq |lhyps| \Rightarrow \exists! d \mid (k, d) \in \text{hist}(c^u_{\phi})$$

Finally, it is worth noticing that the constraint $\ddagger$ is significant only if the considered derivations are in normal form. Therefore, the absence of both the freely accessible identity axiom and the $\sim\sim I$ rule is necessary to the success of our approach.

### 4.3.2 \textit{LGL} grammars & generated language

\textit{LGL} grammars have two parameters, namely a lexicon and an atomic distinguished type $c$. Let $G(\text{lex}, c)$ be a \textit{LGL} grammar and ‘at’ an atomic syntactic type. We say that a sequence of phonetic constants $l=m_1m_2...m_n$ has abstract type ‘at’ within $G$ (i.e. $l \in L_{at}(G)$) iff:

$$\exists x, x_1 \mid x, x_1 \in \text{struct}(m_1, ..., m_n) \land (\vdash (x, x_1) : at; \emptyset)$$

where $\text{struct}(m_1, ..., m_n)$ is the range of phonetic structures built using $\bullet$ operator and whose leaves are $m_1, m_2, ..., m_n$ in that order.

Notice that the convergence of derivations requires the introduction and the simultaneous abstraction of all controlled assumptions related to involved lexical entries.

Finally, checking whether a sequence of phonetic constants $l$ is recognized by the grammar $G$ (i.e. $l \in L(G)$) amounts to verifying that $l$ has abstract type $c$.

### 4.3.3 Example of \textit{LGL} derivations

This section is devoted to the study of a hybrid example ‘More logicians met Godel than physicists knew him’ which involves two complex linguistic phenomena: binding and discontinuity. The analysis of these phenomena within the directional approach constitutes a real challenge for researchers. All proposed solutions are complex insofar as they led to the extension of the core logic either by defining new syntactic connectives (discontinuity connectives: Morrill (2000)) or by introducing additional packages of postulates as in Hendriks (1995). However, our proposal is able to capture such phenomena in an elegant fashion without using any additional material.

Our treatment of binding follows the same ideas of Kayne. R in Kayne (2002) where he argues that the antecedent-pronoun relation (e.g. between \textit{Godel} and \textit{him}) stems from the fact that both enter the derivation together as a doubling constituent ([\textit{Godel}, \textit{him}]) and are subsequently separated after movement. In our system, we account for this idea by defining a linked entry $e_1$ (cf. Fig. 4) associated with the proper noun \textit{Godel}. This entry requires the introduction of two hypotheses (where the first ‘him’ is a pronounced one)
which must be discharged at the same time. Therefore, $e_1$ entry will reach its final position thus making it possible to semantically link the pronoun with its antecedent.

<table>
<thead>
<tr>
<th>Id</th>
<th>$\Phi$-terms</th>
<th>$\Lambda$-terms</th>
<th>Abstract types</th>
<th>Hyps</th>
</tr>
</thead>
<tbody>
<tr>
<td>$e_1$</td>
<td>$A \vdash P_2$, $P_3$ (Godel)</td>
<td>$A \vdash P_2$, $P_3$ (Godel)</td>
<td>($d_{av} \vdash \neg c$) $\vdash \neg \neg c$</td>
<td>$[X \cdot d_{av} \vdash \neg \neg c]$, $[X \cdot d_{av} \vdash \neg \neg c]$</td>
</tr>
<tr>
<td>$e_2$</td>
<td>logicians</td>
<td>Logician</td>
<td>n</td>
<td>()</td>
</tr>
<tr>
<td>$e_3$</td>
<td>physicists</td>
<td>Physicists</td>
<td>n</td>
<td>()</td>
</tr>
<tr>
<td>$e_4$</td>
<td>$Ax \cdot Ay \cdot \Phi (x) \cdot (\Phi (x) \cdot \Phi (y) \cdot \Phi (x) \cdot \Phi (y) \cdot \Phi (x) \cdot \Phi (y))$</td>
<td>$Ax \cdot Ay \cdot \Phi (x) \cdot (\Phi (x) \cdot \Phi (y) \cdot \Phi (x) \cdot \Phi (y) \cdot \Phi (x) \cdot \Phi (y))$</td>
<td>($d_{av} \vdash \neg d_{av}$) $\vdash \neg \neg d_{av}$</td>
<td>$[X \cdot d_{av} \vdash \neg \neg d_{av}]$, $[X \cdot d_{av} \vdash \neg \neg d_{av}]$</td>
</tr>
<tr>
<td>$e_5$</td>
<td>$Ax \cdot Ay \cdot \Phi (x) \cdot (\Phi (x) \cdot \Phi (y) \cdot \Phi (x) \cdot \Phi (y) \cdot \Phi (x) \cdot \Phi (y))$</td>
<td>$Ax \cdot Ay \cdot \Phi (x) \cdot (\Phi (x) \cdot \Phi (y) \cdot \Phi (x) \cdot \Phi (y) \cdot \Phi (x) \cdot \Phi (y))$</td>
<td>($d_{av} \vdash \neg d_{av}$) $\vdash \neg \neg d_{av}$</td>
<td>$[X \cdot d_{av} \vdash \neg \neg d_{av}]$, $[X \cdot d_{av} \vdash \neg \neg d_{av}]$</td>
</tr>
</tbody>
</table>

**FIGURE 4 Example of $LGL$ lexicon**

On the other hand, we capture discontinuity by gathering the different components of a discontinuous expression in the same lexical entry. For instance, entry $e_5$ defines the phonetic and semantic behavior of the discontinuous constituent (more ... than).

We present, in the following, the main steps of our example’s analysis. For the sake of legibility, the bottom-up derivation tree is split into different key parts which will be commented on progressively.

The derivation above starts by introducing the last controlled hypothesis (i.e. the assumption representing the accusative pronoun *him*) of the sequence attached to $e_1$ entry. This hypothesis, then, merges with lexical entry $e_5$ by means of $\neg E$ rule. On the other hand, a partial derivation is built by consecutively combining entry $e_5$ with entries $e_3$ and $e_2$. The resulting sequent then merges with the previous one. The last deduction step does nothing but decorating the type of the introduced hypothesis by a ! marker in order to express its ability to gather with the other controlled hypothesis linked to its proper
entry. At this stage of analysis, only the second controlled hypothesis of \( e_1 \) has been used. Moreover, it was involved in exactly three deduction steps after its introduction, so we can deduce that its current history is: \( \text{hist}(x^{11}) = ((2,3)) \).

\[
\begin{array}{c}
\vdash (\lambda x. \phi \bullet \{ \text{meets} \bullet y \} ) \\
\vdash (\lambda x. \text{Meet}_{\phi}(y, x))
\end{array}
\]

In this second part of analysis, the first controlled assumption linked to \( e_1 \) entry is introduced. Then, it merges with \( e_2 \) entry which represents the past form of the transitive verb *meet*. This branch of the derivation ends by a !L step like the previous one. We can easily check that, at this point of the derivation, the history of \( x^{11} \) assumption is nothing else but \( \text{hist}(x^{11}) = ((1,2)) \).

The partial derivation above stems from merging the two previously presented branches into one tree. Contraction rule is then applied to encapsulate both controlled hypotheses linked to \( e_1 \) in one assumption \( y^{11} \). The current history of this latter compound assumption is: \( \text{hist}(y^{11}) = ((1,4) \cup (2,5)) \).

The whole derivation ends by simultaneously discharging controlled hypotheses linked to entry \( e_1 \) by means of \( \neg \text{IE} \) rule. In fact, the application of this rule is allowed since the side-condition \( \exists \) is entirely verified: as \( y^{11} \)'s history shows, the leftmost hypothesis linked to \( e_1 \) was introduced in the derivation after the rightmost one. The semantic representation of our sentence is computed in tandem. Indeed, the final semantics coincides with the intuitive meaning of the sentence, namely that the set of logicians who met Godel is larger than the range of physicists that knew him.
4.4 Enhancing LGL

It is not difficult to notice that our logic is too flexible as the application of
movement is not constrained. For instance, if we assign the entry below to
the wh-element ‘which’, we can analyze both sentences *which man do you
think the child of *speaks? and ‘which man do you think John loves the child
of _3’*, where the first is ungrammatical.

\[\lambda m \lambda d (\text{which } \bullet_1 n) \bullet_2 \phi(e)\]

\[\Lambda P \Lambda Q \lambda x.P(x) \wedge Q(x)\]

\[n \rightarrow (d_{\text{dat}} \rightarrow c) \rightarrow c \rightarrow d \rightarrow X : d_{\text{dat}} \rightarrow X : d_{\text{dat}}\]

In fact, we need to control displacement operation to rule out extraction from
islands. For that purpose, we propose to enhance LGL with some meta-rules
encoding locality constraints (e.g. SPIC: Specifier Island Condition, SMC:
Shortest Move Condition). We focus in the following on the SPIC defined
in Koopman and Szabolcsi (2000) which stipulates that the moved element
should be a member of the extraction domain (i.e. comp*: transitive closure
of the complement relation, or a specifier of a comp*).

In order to locate the position of the head, the complement and the specifier
inside a phonetic expression, we decorate the building structure connective •
with a mode of composition taken from the set \{<, >\}. This mode points
towards the sub-tree where the head is located: •_< (resp. •_>) if the head is
located on the left (resp. right) sub-tree.

A linked lexical entry which is expected to undergo an overt constituent
movement has a phonetic-term that obeys the following syntax:

\[\lambda x_1 \ldots \lambda x_n \lambda P_\Phi \lambda y_1 \ldots \lambda y_k \cdot g(y_1, \ldots, y_k, f(x_1, \ldots, x_n) \bullet_1 P_\Phi(e))\]

In the expression above, \(x_1, \ldots, x_n, y_1, \ldots, y_k\) (\(n \geq 0, k \geq 0\)) are \(\Phi\)-variables of
arbitrary types, whereas \(P_\Phi\) is a \(\Phi\)-variable of type \(s \rightarrow s\). Moreover, \(f\) (resp. \(g\))
is a function that takes \(n\) (resp. \(k+1\)) \(\Phi\)-terms and builds a phonetic structure
using these parameters together with constants of \(\Sigma\).

Intuitively, this syntax means that our entry will firstly combine with \(n\)
structures \(x_1, \ldots, x_n\) by means of merge operation, thus leading to a maxi-
mal projection \(f(x_1, \ldots, x_n)\). Then, the intermediary sites will be replaced by
traces in the initial structure \(P_\Phi\) and our maximal projection will be placed in
specifier position, hence making it possible to carry out the expected move-
ment. Finally, our resulting constituent can merge with other structures, thus
yielding a complete expression (e.g. whom entry in section 2.2).

Notice that this syntax suits the type specification defined in section 2.2
(points 2 & 3) if we add additional conditions, namely that both types \(t\)
type of intermediary sites) and \(t'\) (type of the D-structure before movement)
are atomic. The first condition (i.e. \(t \in \mathcal{A}\)) follows from constraints proposed
by Koopman and Szabolcsi (Koopman and Szabolcsi (2000)) which forces

\(d_{\text{dat}}\) represent noun phrases with dative case.
moved elements to be maximal projections (i.e. complete expressions). However, the latter condition \(t' \in \mathcal{A}\) is a logical formalization of the merge over move principle Chomsky (1995) which stipulates that merge operation has priority over movement because of its simplicity. Therefore, a structure that will undergo move operation should be complete.

According to the syntax of phonetic terms associated with moved elements, SPIC condition can be encoded in LGL as a pre-condition of \(\neg\text{IE}\) rule (cf. Fig 3) stipulating the inclusion of all occurrences of \(\Phi\)-variable \(c_\Phi\) within the extraction domain of the \(\Phi\)-term \(d_\Phi\). Therefore, adding this meta-rule to LGL prevents us from analyzing the previous ungrammatical sentence.

### 4.5 Conclusion & Future Work

LGL is a new logical formalism which proposes a deductive simulation of Minimalist Program. Our proposal is powerful enough to describe several linguistic phenomena such as medial extraction, binding, ellipsis and discontinuity thanks to using linked lexical entries (related to controlled hypotheses). Moreover, one can solve over-generation problems caused by the freedom of displacement by adding some meta-rules encoding locality constraints.

In addition, it is not difficult to show that these grammars are richer than context free grammars as they are able to generate crossed-dependencies languages (e.g. \(\{a^n b^m c^n d^m \mid n, m \geq 0\}\)). In fact, this latter language is recognized by LGL grammar containing the lexicon below:\(^8\):

\[
\begin{align*}
\vdash & c; p_i (\forall i \in \{1..4\}) \\
\vdash & d_x.d_y.d_z.d_u.x@y@z@u); ty \\
\vdash & d P; d_x.d_y.d_z.d_u. P(x,y,z,u); ty \rightarrow ty \\
\vdash & d P; d_x.d_y.d_z.d_u. P(x,b,y,z,d,u); ty \rightarrow ty
\end{align*}
\]

The next direction to explore concerns the study of LGL formal properties: expressive power, decidability, and complexity. We also intend to build bridges between LGL and other well-known grammatical frameworks (e.g. Minimalist Grammar, TAGs).

Finally, we are developing a meta-linguistic toolkit using Coq proof assistant (Coq Team (2004)), in order to study logical properties of LGL grammars being enhanced with packages of meta-constraints. This toolkit can help users manage complex derivations by automatically handling some technical proofs thanks to powerful computation tools (strategies).

### References


---

\(^8\)In that case, \(\mathcal{A} = \{p_1, p_2, p_3, p_4, c\}\) and \(ty\) denotes the composite type \(p_1 \rightarrow p_2 \rightleftharpoons p_3 \rightarrow p_4 \rightarrow c\)


P-TIME decidability of NL1 with assumptions

MARIA BULIŃSKA

Abstract

Buszkowski (2005) showed that systems of Non-associative Lambek Calculus with finitely many non-logical axioms are decidable in polynomial time and generate context-free languages. The same holds for systems with unary modalities, studied in Moortgat (1997), n-ary operations, and the rule of permutation, studied in Jäger (2004). The polynomial time decidability for Classical Non-associative Lambek Calculus was established by de Groote and Lamarche (2002). We study Non-associative Lambek Calculus with identity enriched with a finite set of assumptions. To prove that this system is decidable in polynomial time we adapt the method used in Buszkowski (2005). The context-freeness of the languages generated of the systems of Non-associative Lambek Calculus is also established.

Keywords

Lambek calculus, P-TIME decidability

5.1 Introduction and preliminaries

Non-logical axioms can be of interest for linguistics for several reasons. We can use them to describe subcategorization in natural language. For instance, restrictive adjectives are a sub-category of adjectives. Further, by enriching Non-associative Lambek Calculus with finitely new axioms, we can improve its expressibility without lacking the nice computational simplicity.

First we describe the formalism of Non-associative Lambek Calculus with identity constant (NL1). Let $\text{At} = \{p, q, r, \ldots\}$ be the denumerable set of atoms (primitive types).

The set of formulas (also called types) $\text{Tp1}$ is defined as the smallest set fulfilling the following conditions:

- $\mathbf{1} \in \text{Tp1}$,
At $\subseteq Tp_1$.

If $A, B \in Tp_1$, then $(A \bullet B) \in Tp_1, (A/B) \in Tp_1, (A\backslash B) \in Tp_1$, where binary connectives $\backslash, /, \bullet$, are called left residuation, right residuation, and product, respectively.

The set of formula structures $STR_1$ is defined recursively as follows:

- $\Lambda \in STR_1$, where $\Lambda$ denotes an empty structure,
- $Tp_1 \subseteq STR_1$; these formula structures are called atomic formula structures,
- if $X, Y \in STR_1$, then $(X \circ Y) \in STR_1$.

We set $(X \circ \Lambda) = (\Lambda \circ X) = X$.

Substructures of a formula structure are defined in the following way:

- $\Lambda$ is the only substructure of $\Lambda$,
- if $X$ is an atomic formula structure, then $\Lambda$ and $X$ are the only substructures of $X$,
- if $X = (X_1 \circ X_2)$, then $X$ and all substructures of $X_1$ and $X_2$ are substructures of $X$.

By $X[Y]$ we denote a formula structure $X$ with a distinguished substructure $Y$, and by $X[Z]$ - the substitution of $Z$ for $Y$ in $X$.

Sequents are formal expressions $X \rightarrow A$ such that $A \in Tp_1$, $X \in STR_1$.

The Gentzen-style axiomatization of the calculus NL1 employs the axiom schemas:

(Id) $A \rightarrow A$   (IR) $\Lambda \rightarrow 1$

and the following rules of inference:

1. $X[\Lambda] \rightarrow A$   $X[1] \rightarrow A$

2. $X[A \circ B] \rightarrow C$   $X[A \bullet B] \rightarrow C$.  (\bulletR) $X \rightarrow A$; $Y \rightarrow B$ $\rightarrow X \circ Y \rightarrow A \bullet B$.

3. $Y \rightarrow A$; $X[B] \rightarrow C$   $X[\circ (A/B)] \rightarrow C$.  (\backslashR) $A \circ X \rightarrow B$ $\rightarrow X \rightarrow A\backslash B$.

4. $X[A] \rightarrow C$; $Y \rightarrow B$   $X[(B/A) \circ Y] \rightarrow C$.  (/R) $X \circ Y \rightarrow A$ $\rightarrow X \rightarrow A/B$.

5. $Y \rightarrow A$; $X[A] \rightarrow B$   $X[Y] \rightarrow B$.  (CUT) $\rightarrow X \rightarrow A$.

For any system $S$ we write $S \vdash X \rightarrow A$ if the sequent $X \rightarrow A$ is derivable in $S$.

The most general models of NL1 are residuated groupoid with identity.
A residuated groupoid with identity is a structure
\[ M = (M, \leq, \cdot, \backslash, /, 1) \]
such that

- \((M, \cdot, 1)\) is a groupoid with identity in which \(a \cdot 1 = a, 1 \cdot a = a\) for all \(a \in M\),
- \((M, \leq)\) is a poset,
- \(\backslash, /\) are binary operations on \(M\) satisfying the equivalences:

\[(RES)\quad ab \leq c \quad \text{iff} \quad b \leq a \backslash c \quad \text{iff} \quad a \leq c / b\]
for all \(a, b, c \in M\).

Every residuated groupoid fulfills the following monotonicity laws:

\[(MON)\quad \text{If} \quad a \leq b \quad \text{then} \quad ca \leq cb \quad \text{and} \quad ac \leq bc\]
\[(MRE)\quad \text{If} \quad a \leq b \quad \text{then} \quad c \backslash a \leq c \backslash b, \quad a / c \leq b / c, \quad b \backslash c \leq a \backslash c, \quad c / b \leq c / c\]
for all \(a, b, c \in M\).

A model is a pair \((M, \mu)\) such that \(M\) is a residuated groupoid with identity and \(\mu\) is an assignment of elements of \(M\) for atoms. One extends \(\mu\) for all formulas:

\[\mu(1) = 1, \quad \mu(A \bullet B) = \mu(A) \cdot \mu(B), \quad \mu(A \backslash B) = \mu(A) \mu(B), \quad \mu(A / B) = \mu(A) / \mu(B).\]

and formula structure:

\[\mu(\Lambda) = \mu(1) = 1, \quad \mu(X \circ Y) = \mu(X) \cdot \mu(Y).\]

A sequent \(X \rightarrow A\) is said to be true in model \((M, \mu)\) if \(\mu(X) \leq \mu(A)\). In particular a sequent \(\Lambda \rightarrow A\) is said to be true in model \((M, \mu)\) if \(1 \leq \mu(A)\).

One can prove the following property for formula structures:

\[\text{(MON−STR)}\quad \text{If} \quad \mu(Y) \leq \mu(Z) \quad \text{then} \quad \mu(X[Y]) \leq \mu(X[Z]).\]

### 5.2 NL1 with assumptions

Let \(\Gamma\) be a set of sequents of the form \(A \rightarrow B\), where \(A, B \in \text{Tp}_1\). By NL1(\(\Gamma\)) we denote the calculus NL1 with additional set \(\Gamma\) of assumptions. NL1 is strongly complete with respect to the residuated groupoids with identity, i.e. all sequents provable in NL1(\(\Gamma\)) are precisely those which are true in all models \((M, \mu)\) in which all sequents from \(\Gamma\) are true. Soundness is easily proved by induction on derivation in NL1(\(\Gamma\)). Completeness follows from the fact that the Lindenbaum algebra of NL1 is a residuated groupoid with identity.

In general, the calculus NL1(\(\Gamma\)) has not the standard sub-formula property, since (CUT) is legal rule in this system. Thus we take into consideration the sub-formula property in some extended form.
Let $T$ be a set of formulas closed under sub-formulas and such that all formulas appearing in $\Gamma$ belong to $T$. By a $T$-sequent we mean a sequent $X \rightarrow A$ such that $A$ and all formulas appearing in $X$ belong to $T$. Now, we can reformulate the sub-formula property as follows:

Every $T$-sequent provable in a system $S$ has a proof in $S$ such that all sequents appearing in this proof are $T$-sequents.

To prove the sub-formula property for $\text{NL1}(\Gamma)$ we will use special models, namely residuated groupoids with identity of cones over given pre-ordered groupoids with identity.

Let $(M, \leq, \cdot)$ be a pre-ordered groupoid, that means, it is a groupoid with a pre-ordering (i.e. a reflexive and transitive relation), satisfying (MON).

A set $P \subseteq M$ is called a cone on $M$ if $a \leq b$ and $b \in P$ entails $a \in P$. Let $C(M)$ denotes the set of cones on $M$.

The operations $\cdot, \setminus, /$ on $C(M)$ are defined as follows:

(M1) $I = \{a \in M : a \leq 1\}$

(M2) $P_1P_2 = \{c \in M : (\exists a \in P_1, b \in P_2) c \leq ab\}$

(M3) $P_1\setminus P_2 = \{c \in M : (\forall a \in P_1) ac \in P_2\}$

(M4) $P_1/P_2 = \{c \in M : (\forall b \in P_2) cb \in P_1\}$.

A structure $(C(M), \subseteq, \cdot, \setminus, /, I)$ is a residuated groupoid with identity. It is called the residuated groupoid with identity of cones over the given pre-ordered groupoid with identity.

Let $\Lambda \in M$ be the set of all formula structures all of whose atomic substructures belong to $T$ and $\Lambda \in M$. If a sequent $X \rightarrow A$ has a proof in $\text{NL1}(\Gamma)$ consisting of $T$-sequents only, we write: $X \rightarrow_T A$.

First, we define on $M$ a relation $\leq_b$. $X \leq_b Y$ denotes $X$ directly reduces to $Y$. The definition of this relation is as follows:

$Y[Z] \leq_b Y[A]$ if $Z \rightarrow_T 1$,

$Y[Z] \leq_b Y[A]$ if $Z \rightarrow_T A$,

$Y[A \bullet B] \leq_b Y[A \circ B]$ if $A \bullet B \in T$.

A pre-ordering $\leq$ on $M$ is defined as a reflexive and transitive closure of the relation $\leq_b$. Then $X \leq Y$ iff there exist $Y_0, \ldots, Y_n, n \geq 0$ such that $X = Y_0, Y = Y_n$ and $Y_{i-1} \leq_b Y_i$, for each $i = 1, \ldots, n$.

Clearly, $(M, \leq, \circ, \Lambda)$ is a pre-ordered groupoid with identity $\Lambda$ fulfilling (MON).

Next, we take into consideration the residuated groupoid of cones with identity $C(M) = (C(M), \subseteq, \cdot, \setminus, /, I)$ over $(M, \leq, \circ, \Lambda)$ defined above. An assignment $\mu$ on $C(M)$ is defined by setting:

$\mu(p) = \{X \in M : X \rightarrow_T p\}$,
for all atoms $p$. One can easily prove that

$$\mu(A) = \{ X \in M : X \rightarrow_X A \},$$

for all $A \in T$.

**Fact 1** Every sequent provable in $\text{NL1}(\Gamma)$ is true in $(C(M), \mu)$.

*Proof.* It suffice to show, that each axiom from $\Gamma$ is true in $(C(M), \mu)$. Assume that $A \rightarrow B$ belongs to $\Gamma$. It yields $A \rightarrow_{\mu} B$. We need to show that $\mu(A) \subseteq \mu(B)$. Let $X \in \mu(A)$. Then, $X \rightarrow_{\mu} A$. By (CUT), we get $X \rightarrow_{\mu} B$, which yields $X \in \mu(B)$. \hfill $\square$

**Lemma 2** The system $\text{NL1}(\Gamma)$ has the extended sub-formula property.

*Proof.* Let $X \rightarrow A$ be a $T$-sequent provable in $\text{NL1}(\Gamma)$. By fact 1 it is true in the model $(C, \mu)$, i.e. $\mu(X) \subseteq \mu(A)$. Since $X \in \mu(X)$, we have $X \in \mu(A)$. But it means $X \rightarrow_{\mu} A$. \hfill $\square$

A sequent is said to be *basic* if it is a $T$-sequent of the form $\Lambda \rightarrow A$, $A \rightarrow B$, $A \circ B \rightarrow C$. Let $\Gamma$ be finite, and let $T$ be a finite set of formulas, closed under sub-formulas and such that $T$ contains all formulas appearing in $\Gamma$. For such $T$ we shall describe an effective procedure which produces all basic sequents derivable in $\text{NL1}(\Gamma)$.

Let $S_0$ consist of all $T$-sequent of the form (Id), all sequents from $\Gamma$, $\Lambda \rightarrow 1$ and all $T$-sequents of the form:

- $1 \circ A \rightarrow A$, $A \circ 1 \rightarrow A$, $A \circ B \rightarrow A \bullet B$, $A \circ (A \setminus B) \rightarrow B$, $(A / B) \circ B \rightarrow A$.

Assume $S_n$ has already been defined. $S_{n+1}$ is $S_n$ enriched with sequents resulting from the following rules:

(S1) if $(A \circ B \rightarrow C) \in S_n$ and $(A \bullet B) \in T$, then $(A \bullet B \rightarrow C) \in S_{n+1}$.

(S2) if $(A \circ X \rightarrow C) \in S_n$ and $(A \setminus C) \in T$, then $(X \rightarrow A \setminus C) \in S_{n+1}$.

(S3) if $(X \circ B \rightarrow C) \in S_n$ and $(C / B) \in T$, then $(X \rightarrow C / B) \in S_{n+1}$.

(S4) if $(A \rightarrow A) \in S_n$ and $(A \circ X \rightarrow C) \in S_n$, then $(X \rightarrow C) \in S_{n+1}$.

(S5) if $(A \rightarrow A) \in S_n$ and $(X \circ A \rightarrow C) \in S_n$, then $(X \rightarrow C) \in S_{n+1}$.

(S6) if $(A \rightarrow B) \in S_n$ and $(B \circ X \rightarrow C) \in S_n$, then $(A \circ X \rightarrow C) \in S_{n+1}$.

(S7) if $(A \rightarrow B) \in S_n$ and $(X \circ B \rightarrow C) \in S_n$, then $(X \circ A \rightarrow C) \in S_{n+1}$.

(S8) if $(A \circ B \rightarrow C) \in S_n$ and $(C \circ D) \in S_n$, then $(A \circ B \rightarrow D) \in S_{n+1}$.

Clearly, $S_n \subseteq S_{n+1}$ for all $n \geq 0$. We define $S_0$ as the join of this chain. $S^T$ is a set of basic sequents, hence it must be finite. It yields $S^T = S_{k+1}$, for the least $k$ such that $S_k = S_{k+1}$, and this $k$ is not greater then the number of basic sequents.

**Fact 3** The set $S^T$ can be constructed in polynomial time.
Proof. Let \( n \) be the cardinality of \( T \). There are \( n, n^2 \) and \( n^3 \) basic sequents of the form \( \Lambda \rightarrow A, A \rightarrow B \) and \( A \circ B \rightarrow C \), respectively. Hence, we have \( m = n^3 + n^2 + n \) basic sequents. The set \( S_0 \) can be constructed in time \( O(n^2) \).

To get \( S_{i+1} \) from \( S_i \) we must check if \((A \circ B \rightarrow C) \in S_i \) and \((A \bullet B) \in T\) which needs at most \( m \) steps for each rule. For example, to close \( S_i \) under \((S1)\) we must check if \((A \circ B \rightarrow C) \in S_i \) and \((A \bullet B) \in T\) which needs at most \( m \) and \( n \) steps, respectively. The sequent \( A \bullet B \rightarrow C \) is added to \( S_{i+1} \) only if it doesn’t belong to this set. To check this fact the next \( m \) steps are needed. The least \( k \) such that \( S^T = S_k \) is at most \( m \). Then finally, we can construct \( S^T \) from \( T \) in time \( O(m^4) = O(n^{12}) \).

By \( S(T) \) we denote the system whose axioms are all sequents from \( S^T \) and whose only inference rule is \((\text{CUT})\). Then, every proof in \( S(T) \) consist of \( T\)-sequents only.

If as premises of \((\text{CUT})\) in the proof in \( S(T) \) of some sequent \( X \rightarrow A \) only sequents without empty antecedents are used, then the length of all sequents in this proof is not greater than the length of \( X \rightarrow A \). But it doesn’t hold if we allow in \((\text{CUT})\) the premises of the form \( \Lambda \rightarrow A \). Therefore we introduce another system \( S'(T) \) whose axioms are all sequents from \( S^T \) and whose only inference rule is \((\text{CUT})\) with premises without empty antecedents, and show the following lemma.

Lemma 4. For any sequent \( X \rightarrow A \), \( S'(T) \vdash X \rightarrow A \) iff \( S(T) \vdash X \rightarrow A \).

Proof. The 'if' direction is evident. To prove the 'only if' direction we show that \( S(T) \) is closed under \((\text{CUT})\), i.e. 

(*) If \( S(T) \vdash X \rightarrow B \) and \( S(T) \vdash Y[B] \rightarrow A \), then \( S(T) \vdash Y[X] \rightarrow A \).

Assume \( S(T) \vdash X \rightarrow B \) and \( S(T) \vdash Y[B] \rightarrow A \).

If \( X \neq \Lambda \), then \( S(T) \vdash Y[X] \rightarrow A \) by definition of \( S(T) \).

If \( X = \Lambda \), then the sequent \( X \rightarrow B \) is of the form \( \Lambda \rightarrow B \) and \( S(T) \vdash \Lambda \rightarrow B \), which means that \( \Lambda \rightarrow B \) is an axiom of \( S(T) \).

To prove (*) proceed by induction on derivation of the second premise: \( Y[B] \rightarrow A \).

If \( Y[B] \rightarrow A \) is an axiom of \( S(T) \), then \( (Y[B] \rightarrow A) \in S^T \). \( S^T \) is closed under \((\text{CUT})\). Hence, \( (Y[\Lambda] \rightarrow A) \in S^T \) which yields \( S(T) \vdash Y[\Lambda] \rightarrow A \).

If \( Y[B] \rightarrow A \) is a conclusion of \((\text{CUT})\) from premises without empty antecedents, then \( Y[B] = Z[Y'] \) and for some \( C \in T \), \( S(T) \vdash Y' \rightarrow C \) and \( S(T) \vdash Z[C] \rightarrow A \). We consider the following cases.

I. \( B \) is contained in \( Y' \). Then \( Y' = Y'[B] \).

1. \( Y'[B] \neq B \). By the induction hypothesis, (*) holds for \( \Lambda \rightarrow B \) and \( Y'[B] \rightarrow C \), so \( S(T) \vdash Y'[\Lambda] \rightarrow C \). Since \( Y'[B] \neq B \), we have \( Y'[\Lambda] \neq \Lambda \). Using \((\text{CUT})\), we get \( S(T) \vdash Z[Y'[\Lambda]] \rightarrow A \), which means \( S(T) \vdash Y[\Lambda] \rightarrow A \).

2. \( Y'[B] = B \). By the induction hypothesis, (*) holds for \( \Lambda \rightarrow B \) and
Lemma 6

If \( S \rightarrow Y \) and \( X \rightarrow A \), we get \( S(T)^- \vdash Z[A] \rightarrow A \), which means \( S(T)^- \vdash Y[A] \rightarrow A \).

II. \( B \) and \( Y' \) do not overlap. Then \( B \) is contained in \( Z \) and does not overlap \( C \) in \( Z \). We write \( Z[C] = Z[B, C] \). From the assumption we have \( Y' \neq \Lambda \). By induction hypothesis, (*) holds for \( \Lambda \rightarrow B \) and \( Z[B, C] \rightarrow A \), so \( S(T)^- \vdash Z[A, C] \rightarrow A \). By (CUT), \( S(T)^- \vdash Z[A, Y'] \rightarrow A \), which means \( S(T)^- \vdash Y'[A] \rightarrow A \).

\[ \square \]

Corollary 5

Every basic sequent provable in \( S(T) \) belongs to \( S^T \).

Proof. We proceed by induction on proofs in \( S(T) \). Assume \( X \rightarrow A \) is a basic sequent derivable in \( S(T) \). If \( X \rightarrow A \) is an axiom of \( S(T) \), then \( (X \rightarrow A) \in S^T \).

If \( X \rightarrow A \) is a conclusion of (CUT), we consider three cases.

1. \( X = \Lambda \). By lemma 4, \( \Lambda \rightarrow A \) has a proof in \( S(T)^- \). Hence \( \Lambda \rightarrow A \) is an axiom, which means \( (\Lambda \rightarrow A) \in S^T \).

2. \( X = B \). By lemma 4, there exists a proof such that \( B \rightarrow A \) is a conclusion from premises \( B \rightarrow C \) and \( C \rightarrow A \), where \( C \neq \Lambda \). Since proofs in \( S(T) \) consist with \( T \)-sequents only, \( B \rightarrow C \) and \( C \rightarrow A \) are basic sequents. By induction hypothesis, \( (B \rightarrow C) \in S^T \) and \( (C \rightarrow A) \in S^T \). \( S^T \) is closed under (CUT), so \( (B \rightarrow A) \in S^T \).

3. \( X = B \circ C \). By lemma 4, there exists a proof such that \( B \circ C \rightarrow A \) is a consequence from premises without empty premises. Hence, they are of the form: \( (B \circ C \rightarrow D, D \rightarrow A) \) or \( (B \rightarrow D, D \circ C \rightarrow A) \) or \( (C \rightarrow D, B \circ D \rightarrow A) \). By the same argument as in (2), in each case, we get \( (B \circ C \rightarrow A) \in S^T \).

Now, we can state an interpolation lemma for \( S(T) \).

Lemma 6

If \( S(T) \vdash X[Y] \rightarrow A \), then there exists \( D \in T \) such that \( S(T) \vdash Y \rightarrow D \) and \( S(T) \vdash X[D] \rightarrow A \).

Proof. We proceed by induction on proofs in \( S(T) \).

I. Assume \( X[Y] \rightarrow A \) is an axiom of \( S(T) \). We consider the following cases.

1. \( X[Y] = Y \). Then \( Y = X \) (observe, that this case includes sub case \( X = \Lambda \)). We set \( D = A \). We have \( S(T) \vdash X \rightarrow A \) from assumption and \( S(T) \vdash A \rightarrow A \), since \( (A \rightarrow A) \in S^T \).

2. \( X[Y] = B, Y = \Lambda \). Then \( X[Y] = X[\Lambda] = B = B \circ \Lambda \) or \( X[Y] = \Lambda \circ B \) and \( D = 1 \). We have \( S(T) \vdash \Lambda \rightarrow 1 \) and \( S(T) \vdash B \rightarrow A \). (\( B \circ 1 \rightarrow B \)) \in \( S^T \), so \( S(T) \vdash B \circ 1 \rightarrow B \). Using (CUT) we get \( S(T) \vdash X[1] \rightarrow A \). For \( X[Y] = \Lambda \circ B \) the argument is dual.
(3) $X[Y] = B \circ C, Y \neq \Lambda$. Then $Y = B$ or $Y = C$, hence $D = B$ or $D = C$, respectively.

(4) $X[Y] = B \circ C, Y = \Lambda$. Then $X[\Lambda]$ has one of the form: $\Lambda \circ (B \circ C)$, $(B \circ C) \circ \Lambda, (\Lambda \circ B) \circ C, (B \circ \Lambda \circ C), B \circ (\Lambda \circ C), B \circ (C \circ A)$. In all these cases we set $D = F$. For example, if $X[\Lambda] = \Lambda \circ (B \circ C)$, we have $S(T) \vdash \Lambda \rightarrow 1$ and using (CUT) to $S(T) \vdash B \circ C \rightarrow A$ and $S(T) \vdash 1 \circ A \rightarrow A$, we get $S(T) \vdash 1 \circ (B \circ C) \rightarrow A$.

II. Assume $X[Y] \rightarrow A$ is a conclusion of (CUT). Then $X[Y] = Z[Y']$ and for some $B \in T$: $S(T) \vdash Y' \rightarrow B$ and $S(T) \vdash Z[B] \rightarrow A$.

In this part the proof is analogous to the proof of lemma 2 in Buszkowski (2005). The following cases are considered.

(1) $Y$ is contained in $Y'$. Then $Y' = Y'[Y]$. By the induction hypothesis, there exists $D \in T$ such that $S(T) \vdash Y \rightarrow D$ and $S(T) \vdash Y'[D] \rightarrow B$. Using (CUT) with the premises $Z[B] \rightarrow A$ and $Y'[D] \rightarrow B$ we get $S(T) \vdash Z[Y'[D]] \rightarrow A$, which means $S(T) \vdash X[D] \rightarrow A$.

(2) $Y'$ is contained in $Y$. Then $X[Y] = X[Y[Y']] = Z[Y']$ and $Z[B] = X[Y[B]]$. By the induction hypothesis, there exists $D \in T$ such that $S(T) \vdash Y[B] \rightarrow D$ and $S(T) \vdash X[D] \rightarrow A$. Using (CUT) with the premises $Y' \rightarrow B$ and $Y[B] \rightarrow D$ we get $S(T) \vdash Y[Y'] \rightarrow D$.

(3) $Y$ and $Y'$ do not overlap. Then $Y$ is contained in $Z$ and does not overlap $B$ in $Z$. We write $Z[B] = Z[B, Y]$. By the induction hypothesis, there exists $D \in T$ such that $S(T) \vdash Y \rightarrow D$ and $S(T) \vdash Z[B, D] \rightarrow A$. Using (CUT) with the premises $Y' \rightarrow B$ and $Z[B, D] \rightarrow B$ we get $S(T) \vdash Z[Y', D] \rightarrow A$, which means $S(T) \vdash X[D] \rightarrow A$.

\[ \square \]

**Lemma 7** For any $T$-sequent $X \rightarrow A$, $X \rightarrow T \rightarrow A$ iff $S(T) \vdash X \rightarrow A$.

**Proof.** Recall, that $X \rightarrow T \rightarrow A$ means that the sequent $X \rightarrow A$ has the proof in $NL_1(\Gamma)$ consisting with $T$-sequents only.

To prove "if" direction observe that $X \rightarrow T \rightarrow A$, for all sequents $X \rightarrow A$ in $S^T$.

The $T$-sequents which are axioms of $NL_1(\Gamma)$ belong to $S_0$. Thus, to prove the "only if" direction it suffices to show that all inference rules of $NL_1(\Gamma)$, restricted to $T$-sequents, are admissible in $S(T)$. For example, let us consider (IL). Assume $X[\Lambda] \rightarrow A$. By lemma 6, there exist $D \in T$ such that $S(T) \vdash \Lambda \rightarrow D$ and $S(T) \vdash X[D] \rightarrow A$. Since $(D \circ 1 \rightarrow D) \in S^T$, then $S(T) \vdash D \circ 1 \rightarrow D$. By two applications of (CUT), we get $S(T) \vdash X[\Lambda \circ 1] \rightarrow A$, which means $S(T) \vdash X[1] \rightarrow A$.

\[ \square \]

**Theorem 8** If $\Gamma$ is finite, then $NL_1(\Gamma)$ is decidable in polynomial time.

**Proof.** Let $\Gamma$ be a finite set of sequents of the form $B \rightarrow C$ and let $X \rightarrow A$ be a sequent. Let $n$ be the number of logical constants and atoms in $X \rightarrow A$
and $\Gamma$. As $T$ we choose the set of all sub-formulas of formulas appearing in $X \to A$ and formulas appearing in $\Gamma$. Since the number of sub-formulas of any formula $B$ is equal to the number of logical constants and atoms in $B$, $T$ has $n$ elements and we can construct it in time $O(n^2)$. By lemma 2, $\text{NL1}(\Gamma) \vdash X \to A$ iff $X \to_T A$. By lemma 7, $X \to_T A$ iff $S(T) \vdash X \to A$. Proofs in $S(T)$ are actually derivation trees of a context-free grammar whose production rules are the reversed sequents from $S^T$. Checking derivability in context-free grammars is P-TIME decidable. For example, by known CYK algorithm, it can be done in time not exceed $k \cdot n^3$, where $k$ is the size of $S^T$. By the proof of fact 3, the size of $S^T$ is at most $O(n^3)$ and $S^T$ can be constructed in $O(n^{12})$. Hence, the total time is $O(n^{12})$, i.e. $\text{NL1}(\Gamma)$ is P-TIME decidable.

By theorem 8, we have immediately that languages generated by the categorial grammar based on the system $\text{NL1}(\Gamma)$ are context-free. In Buszkowski (2005) the analogous result was established for $\text{NL}(\Gamma)$, $\text{NL}(\Gamma)$ with permutation rule and Generalized Lambek Calculus ($\text{GLC}(\Gamma)$). The context-freeness of the languages generated by Non-associative Lambek Calculus were studied by Buszkowski (1986), Kandulski (1988) and Jäger (2004). Bulińska (2005) obtained the weak equivalence of context-free grammars and grammars based on the associative Lambek calculus with finite set of simple non-logical axioms of the form $p \to q$, where $p, q$ are primitive types.

References


Program transformations for optimization of parsing algorithms and other weighted logic programs

JASON EISNER AND JOHN BLATZ

Abstract

Dynamic programming algorithms in statistical natural language processing can be easily described as weighted logic programs. We give a notation and semantics for such programs. We then describe several source-to-source transformations that affect a program’s efficiency, primarily by rearranging computations for better reuse or by changing the search strategy. We present practical examples of using these transformations, mainly to optimize context-free parsing algorithms, and we formalize them for use with new weighted logic programs.

Specifically, we define weighted versions of the folding and unfolding transformations, whose unweighted versions are used in the logic programming and deductive database communities. We then present a novel transformation called speculation—a powerful generalization of folding that is motivated by gap-passing in categorial grammar. Finally, we give a simpler and more powerful formulation of the magic templates transformation.¹

Keywords  WEIGHTED LOGIC PROGRAMMING, DYNAMIC PROGRAMMING, PROGRAM TRANSFORMATION, PARSING ALGORITHMS

6.1 Introduction

In this paper, we show how some algorithmic efficiency tricks used in the natural language processing (NLP) community, particularly for parsing, can be regarded as specific instances of transformations on weighted logic programs.

¹This material is based upon work supported by the National Science Foundation under Grants No. 0313193 and 0347822 to the first author.
We define weighted logic programs and sketch the general form of the transformations, enabling their application to new programs in NLP and other domains. Several of the transformations (folding, unfolding, magic templates) have been known in the logic programming community, but are generalized here to our weighted framework and applied to NLP algorithms. We also present a powerful generalization of folding—speculation—which appears new and is able to derive some important parsing algorithms.

We also formalize these transformations in a way that we find more intuitive than conventional presentations. Influenced by the mechanisms of categorial grammar, we introduce “slashed” terms whose values may be regarded as functions. These slashed terms greatly simplify our constructions. In general, our work can be connected to the well-established literature on grammar transformation.

The framework that we use for specifying the weighted logic programs is roughly based on that of Dyna (Eisner et al., 2005), an implemented system that can compile such specifications into efficient C++. Some of the programs could also be handled by PRISM (Zhou and Sato, 2003), an implemented probabilistic Prolog.

It is especially useful to have general optimization techniques for dynamic programming algorithms (a special case in our framework), because computational linguists regularly propose new such algorithms. Dynamic programming is used to parse many different grammar formalisms, and in syntax-based approaches to machine translation and language modeling. It is also used in finite-state methods, stack decoding, and grammar induction.

One might select program transformations either manually or automatically. Our goal here is simply to illustrate the search space of semantically equivalent programs. We do not address the practical question of searching this space—that is, the question of where and when to apply the transformations. For some programs and their typical inputs, a transformation will speed a program up (at least on some machines); in other cases, it will slow it down. The actual effect can of course be determined empirically by running the transformed program on typical inputs (or, in some cases, can be reasonably well predicted from runtime profiles of the untransformed program). Thus, one could in principle use automated methods, such as stochastic local search, to search for sets of transformations that provide good practical speedups.

### 6.2 Weighted Logic Programming

Before moving to the actual transformations, we will take several pages to describe our proposed formalism of weighted logic programming.
6.2.1 Logical Specification of Dynamic Programs

We will use context-free parsing as a simple running example. Recall that one can write a logic program for CKY recognition (Younger, 1967) as follows, where \texttt{constit}(X,I,K) is provable iff the context-free grammar (CFG), starting at nonterminal X, can generate the input substring from position I to position K. The capitalized symbols are variables.

\begin{verbatim}
constit(X,I,K) :- rewrite(X,W), word(W,I,K).
constit(X,I,K) :- rewrite(X,Y,Z), constit(Y,I,J), constit(Z,J,K).
goal :- constit(s,0,N), length(N).

rewrite(s,np,vp). % tiny grammar
rewrite(np,det,n).
rewrite(np,“Dumbo”).
rewrite(np,”flies”).
rewrite(vp,”flies”).

word(“Dumbo”,0,1). % tiny input sentence
word(“flies”,1,2).
length(2).
\end{verbatim}

For example, the second line permits us to prove the proposition \texttt{constit}(X,I,K) once we can prove that there exist constituents \texttt{constit}(Y,I,J) and \texttt{constit}(Z,J,K)—which are adjacent—as well as a context-free grammar rule \texttt{rewrite}(X,Y,Z) (i.e, X \rightarrow Y Z) to combine them. This deduction is permitted for any specific values of X,Y,Z (presumably nonterminals of the grammar) and I,J,K (presumably positions in the sentence).

We suppose in this paper that the whole program above is specified at compile time. In practice, one might instead wait until runtime to provide the description of the sentence (the word and length facts) and perhaps even of the grammar (the rewrite facts). In this case our transformations could be used only on the part of the program specified at compile time.\(^3\)

The basic objects of the program are terms, defined in the usual way (as in Prolog). Following parsing terminology, we refer to some terms as items; these are terms that the program might prove in the course of its execution, such as \texttt{constit(np,0,1)} but not \texttt{np} (which appears only as a sub-term of items) nor \texttt{constit(foo(bar),baz)}.\(^4\) Each line in the program is an inference rule (or

\(^2\)By convention, we regard positions as falling between input words, so that the substring from I to J is immediately adjacent to the substring from J to K.

\(^3\)This is generally safe provided that the runtime rules may not define in the head, nor evaluate in the body, any term that unifies with the head of a compile-time rule. It is common to assume further that all runtime rules are facts, known collectively as the database.

\(^4\)It is of course impossible to determine precisely which terms the program will prove without running it. It is merely helpful to refer to terms as items when we are discussing their provability or, in the case of weighted logic programs, their value. (“Item” does have a more formal meaning...
Each of the inference rules in the above example is range-restricted. In the jargon of logic programming, this means that all variables (capitalized) in the rule’s left-hand side (rule head) also appear in its right-hand side (rule body). A rule with an empty body is called a fact. If all rules are range-restricted, then all provable terms are ground terms, i.e., terms such as const(s,0,2) that do not contain any variables.

Logic programs restricted in this way correspond to the “grammatical deduction systems” discussed by Shieber et al. (1995). Sikkel (1997) gives many parsing algorithms in this form. More generally, programs that consist entirely of range-restricted rules correspond to conventional dynamic programming algorithms, and we may refer to them informally as dynamic programs.

Dynamic programs can be evaluated by various techniques. The specific technique chosen is not of concern to this paper except in section 6.6. However, for most NLP algorithms, it is common to use a bottom-up or forward chaining strategy such as the one given by Shieber et al., which iteratively proves all transitive consequences of the facts in the program. In the example above, forward chaining starts with the word, rewrite, and length facts and derives successively wider constit items, eventually deriving goal iff the input sentence is grammatical. This corresponds to chart parsing, with the role of the chart being played by a data structure that remembers which items have been proved so far.5

This paper deals with general logic programs, not just dynamic programs. For example, one may wish to state once and for all that an “epsilon” word is available at every position K in the sentence: word(epsilon,K,K). We allow this because it will be convenient for most of our transformations to introduce new non-range-restricted rules, which can derive non-ground items such as word(epsilon,K,K). The above execution strategies continue to apply, but the presence of non-ground items means that they must now use unification matching to find previously derived terms of a given form. For example, if the non-ground item word(epsilon,K,K) has been derived, representing an infinite collection of ground terms, then if the program looks up the set of terms in the chart matching word(W,2,K), it should find (at least) word(epsilon,2,2).

5An alternative strategy is Prolog’s top-down backward-chaining strategy, which starts by trying to prove goal and tries to prove other items as subgoals. However, this strategy will waste exponential time by re-deriving the same constituents in different contexts, or will fail to terminate if the grammar is left-recursive. It may be rescued by memoization, also known as “tabling,” which re-introduces a chart (Sagonas et al., 1994).
One can often eliminate non-range-restricted rules (in particular, the ones we introduce) to obtain a semantically equivalent dynamic program, but we do not here explore transformations for doing so systematically.

### 6.2.2 Weighted Logic Programs

We now define our notion of weighted logic programs, of which the most useful in NLP are the semiring-weighted dynamic programs discussed by Goodman (1999) and Eisner et al. (2005). See the latter paper for a discussion of relevant work on deductive databases with aggregation (e.g., Fitting, 2002, Van Gelder, 1992, Ross and Sagiv, 1992).

In a weighted logic program, each provable item has a value. Our running example is the inside algorithm for context-free parsing:

- **constit(X,I,K)** += rewrite(X,W) * word(W,I,K).
- goal += constit(s,0,N) * length(N).

```plaintext
rewrite(s,np,vp) = 1. % p(s → np vp | s)
rewrite(np,deny,n) = 0.5. % p(np → deny n | np)
rewrite(np,"Dumbo") = 0.4. % p(np → "Dumbo" | np)
rewrite(np,"flies") = 0.1. % p(vp → "flies" | vp)
rewrite(vp,"flies") = 1. % p(vp → "flies" | vp)
word("Dumbo",0,1) = 1. % 1 for all words in the sentence
word("flies",1,2) = 1.
length(2) = 1.
```

This looks just like the unweighted logic program in section 6.2.1, except that now the body of each inference rule is an arbitrary expression, and the :- operator is replaced by an aggregation operator such as + or max. One might call these rules “Horn equations,” by analogy with the (definite) Horn clauses of the unweighted case. A fact is now a rule whose body is a constant or an expression on constants.

To understand the meaning of the above program, consider for example the item constit(s,0,2). The old version of line 2 allowed one to prove constit(s,0,2) if rewrite(s,Y,Z), constit(Y,0,J), and constit(Z,J,2) were all true for at least one triple Y,Z,J. The new version of line 2 instead defines the value of constit(s,0,2)—or more precisely, as

\[ \sum_{Y,Z,J} \text{rewrite}(s,Y,Z) \times \text{constit}(Y,0,J) \times \text{constit}(Z,J,2) \]

The aggregation operator + requires a sum over all ways of grounding the variables that appear only in the rule body, namely Y, Z, and J. The rest of the value of constit(s,0,2) is added in by line 1. We will formalize all of this in section 6.2.3 below.
To put this another way, one way of grounding line 2 (i.e., one way of substituting a ground term for each of its variables) is \( \text{constit}(s,0,2) \xrightarrow{\text{rewrite}} (s,\text{np},\text{vp}) \times \text{constit}(\text{np},0,1) \times \text{constit}(\text{vp},1,2) \). Therefore, one operand to \( \xrightarrow{+} \) in defining the value of \( \text{constit}(s,0,2) \) will be the value (if defined) of \( \text{rewrite}(s,\text{np},\text{vp}) \times \text{constit}(\text{np},0,1) \times \text{constit}(\text{vp},1,2) \).

The result—for this program—is that the computed value of \( \text{constit}(X,I,J) \) will be the traditional inside probability \( \beta_k(I,J) \) for a particular input sentence and grammar.\(^6\)

If the heads of two rules unify, then the rules must use the same aggregation operator, to guarantee that each provable term’s value is aggregated in a consistent way. Each \( \text{constit}(...) \) term above is aggregated with \( \xrightarrow{+} \).

Substituting \( \text{max}= \) for \( \xrightarrow{+} \) throughout the program would find Viterbi probabilities (best derivation) rather than inside probabilities (sum over derivations). Similarly, we can obtain the unweighted recognizer of section 6.2.1 by writing expressions over boolean values:\(^7\)

\[
\text{constit}(X,I,K) \xrightarrow{=} \text{rewrite}(X,Y,Z) \land \text{constit}(Y,I,J) \land \text{constit}(Z,J,K).
\]

Of course, these programs are all essentially recognizers rather than parsers. They only compute a boolean or real value for goal. To recover actual parse trees, one can extract the proof trees of goal. To make the parse trees accessible to the program itself, one can define a separate item \( \text{parse}(\text{constit}(X,I,K)) \) whose value is a tree.\(^8\) We do not give the details here to avoid introducing new notation and issues that are orthogonal to the scope of this paper.

The above examples, like most of our examples, can be handled by the framework of Goodman (1999). However, we allow a wider class of rules. Goodman allows only range-restricted rules (cf. our section 6.2.1), and he requires all values to fall in a single semiring and all rules to use only the semiring operations. The latter requirements—in particular the distributivity property of the semiring—imply that an item’s value can be found by separately computing values for all of its complete proof trees and then aggregating them at the end. That is not the case for neural networks, game trees,

---

\(^6\) However, unlike probabilistic programming languages (Zhou and Sato, 2003), we do not enforce that values be reals in \([0,1]\) or have probabilistic interpretations.

\(^7\) Using \( | \) for “or” and \( \land \) for “and.” The aggregation operators \( \xrightarrow{=} \) and \( \xrightarrow{\land} \) can be regarded as implementing existential and universal quantification.

\(^8\) Another option is to say that the value of \( \text{constit}(X,I,K) \) is not just a number but a \( (\text{number},\text{tree}) \) pair, and to define \( \text{max}= \) over such pairs Goodman (1999). This resembles the use of semantic attachments to build output in programming language parsers. However, it requires building a tree (indeed, many trees, of which the best is kept) for each \( \text{constit} \), including constituents that do not appear in the final parse. Our preferred scheme is to hold the best tree in a separate \( \text{parse}(\text{constit}(X,I,K)) \) item. Then we can choose to use backward chaining, or the magic templates transformation of section 6.6, to limit our computation of parse trees to those that are needed to assemble the final tree, \( \text{parse}(\text{goal}) \).
practical NLP systems that mix summation and maximization, or other useful systems of equations that can be handled in our more general framework.

6.2.3 Semantics of Weighted Logic Programs

We now formalize the semantics of a weighted logic program, and define what it means for a program transformation to preserve the semantics. Readers who are interested in the actual transformations may skip this section, except for the brief definitions of the special aggregation operator = and of side conditions.

In an unweighted logic program, the semantics is the set of provable ground terms. For a weighted logic program, the semantics is a partial function, the valuation function, that maps each provable ground term \( r \) to a value \([r]\). All items in our example above take values in \( \mathbb{R} \). However, one could use values of any type or of multiple types.

The domain of the valuation function \([\cdot]\) is the set of ground terms for which there exist finite proofs under the unweighted version of the program. We extend \([\cdot]\) in the obvious way to expressions on provable ground terms: for example, \([x * y] \equiv [x] * [y]\) provided that \([x]\) and \([y]\) are defined.

For each ground term \( r \) that is provable in program \( P \), let \( P(r) \) be the non-empty multiset of all expressions \( E \), over provable ground terms, such that \( r \oplus r = E \) grounds some rule of \( P \). Here \( \oplus \) denotes the single aggregation operator shared by all those rules.

We now interpret the weighted rules as a set of simultaneous equations that constrain the \([\cdot]\) function. If \( \oplus = + \), then we require that

\[
[r] = \sum_{E \in P(r)} [E]
\]

(perhaps permitting \([r] = \infty \) if the sum diverges). More generally, we require that

\[
[r] = [E_1] \oplus_r [E_2] \oplus_r \ldots
\]

where \( P(r) = \{E_1, E_2, \ldots\} \). For this to be well-defined, \( \oplus \) must be associative and commutative. If \( \oplus = \) is the special operator =, as in the final rules of our example, then we set \([r] = [E_1] \) if \( P(r) \) is a singleton set \( \{E_1\} \), and generate a runtime error otherwise.

Example. In the example of section 6.2.2, lines 1–2, this means that for any particular \( X, I, K \) for which \( \text{constit}(X, I, K) \) is a provable item, \([\text{constit}(X, I, K)]\)

\[\text{equals}\]

\[\text{Note that if a non-ground term can be proved under the program, so can any one of the infinitely many ground terms that instantiates (specializes) that non-ground term. Our formal semantics are described in terms of these ground terms only.}\]
\[
\sum_w \llbracket \text{rewrite}(X,W) \rrbracket \ast \llbracket \text{word}(W,J,J) \rrbracket \\
+ \sum_{X,Y,Z} \llbracket \text{rewrite}(X,Y,Z) \rrbracket \ast \llbracket \text{constit}(Y,J,J) \rrbracket \ast \llbracket \text{constit}(Z,J,K) \rrbracket
\]

where, for example, the second summation ranges over term triples \(J, Y, Z\) such that the summand has a value. We sum over \(J, Y, Z\) because they do not appear in the rule’s head \(\text{constit}(X, I, K)\), which is being defined.

**Remark.** Our constraints on the valuation function \(\llbracket \cdot \rrbracket\) are equivalent to saying that it is a fixed point of an “equational update” operator \(T_P\),\(^{10}\) which acts on valuation functions \(I\) and is analogous to the “monotone consequence” operator for unweighted logic programs. Such a fixed point need not be unique.\(^{11}\) Operationally, one may seek a single fixpoint by initializing \(I = \{\}\), repeatedly updating \(I\) to \(T_P(I)\), and hoping for convergence. That is the basic idea of the forward-chaining algorithm in section 6.2.4 below.

**Side conditions.** A mechanism for handling “side conditions” (e.g., Goodman, 1999) is to use rules like\(^{12}\)

\[
a\mathrel{+}= b \ast c \text{ whenever } \emptyset d.
\]

We define \(\llbracket b \ast c \text{ whenever } \emptyset d \rrbracket \overset{\text{def}}{=} \llbracket b \ast c \rrbracket\), independent of the value of \(d\). But by our earlier definitions, it will appear in \(P(a)\) and be added into \(\llbracket a \rrbracket\) only if the side condition \(d\), along with \(b\) and \(c\), is provable.

**Definition.** Roughly speaking, a program transformation \(P \rightarrow P'\) is said to be **semantics-preserving** iff it preserves the partial function \(\llbracket \cdot \rrbracket\). In other words, exactly the same ground terms must be provable under both programs, and they must have the same values.

We make two adjustments to this rough definition. First, for generality, we must handle the case where \(P\) and \(P'\) do not both have uniquely determined semantics. In general, we say that the transformation is semantics-preserving iff it preserves the set of valuation functions.

Second, we would like a program transformation to be able to introduce new provable items for its own use. Therefore, we only require that it preserve

\(^{10}\)That is, \(\llbracket \cdot \rrbracket = T_P(\llbracket \cdot \rrbracket)\). Given a valuation function \(I\), \(T_P(I)\) is defined as follows: for ordinary ground terms \(r\), put

\[
(T_P(I))(r) = \sum_{E \text{ such that } E \text{ is a ground expression where} \llbracket E \rrbracket \text{ is defined and } r \ominus r = E \text{ grounds some rule of } P} I(E)
\]

if this sum is non-empty, and leave it undefined otherwise. Then extend \(T_P(I)\) over expressions as usual.

\(^{11}\)There is a rich line of research that attempts to more precisely characterize which fixed point gives the “intuitive” semantics of a logic program with negation or aggregation (see e.g. Fitting, 2002, Van Gelder, 1992, Ross and Sagiv, 1992).

\(^{12}\)whenever \(\emptyset d\) is defined to mean “whenever \(d\) is provable,” whereas \(\text{whenever } d\) would mean “whenever \(d\)’s value is true.” The latter construction is also useful, but not needed in this paper.
the restriction of $[\cdot]$ to the Herbrand base of $\mathcal{P}$ (more precisely, to the Herbrand base of all expressible constants and the functors in $\mathcal{P}$). Thus, a transformed version of the inside algorithm would be allowed to prove additional $\text{temp}(\ldots)$ items, but not additional $\text{constit}(\ldots)$ items. The user may therefore safely interrogate the transformed program to find out whether $\text{constit}(np,0.5)$ is provable and if so, what its value is.

Notice that a two-step transformation $\mathcal{P} \rightarrow \mathcal{P}'' \rightarrow \mathcal{P}'$ might introduce new $\text{temp}(\ldots)$ items in the first step and eliminate them in the second. This composite transformation may still be semantics preserving even though its second step $\mathcal{P}'' \rightarrow \mathcal{P}'$ is not.

All of the transformations developed in this paper are intended to be semantics-preserving (except for rule elimination and magic templates, which preserve the semantics of only a subset of the ground terms). To prove this formally, one would show that every fixed point of $T_{\mathcal{P}'}$ is also a fixed point of $T_{\mathcal{P}}$, when restricted to the Herbrand base of $\mathcal{P}$, and conversely, that every fixed point of $\mathcal{P}$ can be extended to a fixed point of $T_{\mathcal{P}'}$.

### 6.2.4 Computing Semantics by Forward Chaining

A basic strategy for computing a semantic interpretation is “forward chaining.” The idea is to maintain current values for all proved items, and to propagate updates to these values, from the body of a rule to its head, until all the equations are satisfied. This may be done in any order, or in parallel (as for the equational update operator of section 6.2.3). Note that in the presence of cycles such as $x += 0.9 \times x$, the process can still converge numerically in finite time (to finite values or to $\infty$, representing divergence). Indeed, the forward chaining strategy terminates in practice for many programs of practical interest.\(^{13}\)

As already noted in section 6.2.1, Shieber et al. (1995) gave a forward chaining algorithm (elsewhere called “semi-naive bottom-up evaluation”) for unweighted dynamic programs. Eisner et al. (2005) extended this to handle arbitrary semiring-weighted dynamic programs. Goodman (1999) gave a mixed algorithm.

Dealing with our full class of weighted logic programs—not just semiring-weighted dynamic programs—is a substantial generalization. Once we allow inference rules that are not range-restricted, the algorithm must derive non-ground items and store them and their values in the chart, and obtain the value of $\text{foo}(3,3)$, if not explicitly derived, by “backing off” to the derived value of non-ground items such as $\text{foo}(X,X)$ or $\text{foo}(X,3)$, which are preferred in turn to

\(^{13}\)Of course, no strategy can possibly terminate on all programs, because the language (even when limited to unweighted range-restricted rules) is powerful enough to construct arbitrary Turing machines. We remark that forward chaining may fail to terminate either because of oscillation or because infinitely many items are derived (e.g., $s(N) = N$).
the less specific foo(X,Y). Once we drop the restriction to semirings, the algorithm must propagate arbitrary updates (notice that it is not trivial to update the result of max= if one of its operands decreases). Certain aggregation operators also allow important optimizations thanks to their special properties such as distributivity and idempotence. Finally, we may wish to allow rules such as reciprocal(X) = 1/X that cannot be handled at all by forward chaining. We defer all these algorithmic details to a separate paper, focusing instead on the denotational semantics.

6.3 Folding: Factoring Out Subexpressions

Weighted logic programs are schemata that define possibly infinite systems of simultaneous equations. Finite systems of equations can often be rearranged without affecting their solutions (e.g., Gaussian elimination). In the same way, weighted logic programs can be transformed to obtain new programs with better runtime.

Notation. We will henceforth adopt a convention of underlining any variables that appear only in a rule’s body, to more clearly indicate the range of the summation. We will also underline any variables that appear only in the rule’s head; these indicate that the rule is not range-restricted.

Example. Consider first our previous rule from section 6.2.2,
\[ \text{constit}(X, I, K) \leftarrow \text{rewrite}(X, Y, Z) \cdot \text{constit}(Y, I, J) \cdot \text{constit}(Z, J, K). \]

If the grammar has \( N \) nonterminals, and the input is an \( n \)-word sentence or an \( n \)-state lattice, then the above rule can be grounded in only \( O(N^3 \cdot n^3) \) different ways. For this—and the other parsing programs we consider here—it turns out that the runtime of forward chaining can be kept down to \( O(1) \) time per grounding.\(^{14}\) Thus the runtime is \( O(N^3 \cdot n^3) \).

However, the following pair of rules is equivalent:
\[ \text{temp}(X, Y, Z, I, J) = \text{rewrite}(X, Y, Z) \cdot \text{constit}(Y, I, J). \]
\[ \text{constit}(X, I, K) \leftarrow \text{temp}(X, Y, Z, I, J) \cdot \text{constit}(Z, J, K). \]

We have just performed a weighted version of the classical folding transformation for logic programs (Tamaki and Sato, 1984). The original body expression would be explicitly parenthesized as \( \text{constit}(X, I, K) \leftarrow \text{rewrite}(X, Y, Z) \cdot \text{constit}(Y, I, J) \cdot \text{constit}(Z, J, K) \); we have simply introduced a “temporary item” (like a temporary variable in a traditional language) to hold the result of the parenthesized subexpression, then “folded” that temporary item into the computation.

\(^{14}\)Assuming that the grammar is acyclic (in that it has no unary rule cycles) and so is the input lattice. Even without such assumptions, a meta-theorem of McAllester (1999) allows one to derive asymptotic run-times of appropriately-indexed forward chaining from the number of instantiations. However, that meta-theorem applies only to unweighted dynamic. Similar results in the weighted case require acyclicity. Then one can use the two-phase method of Goodman (1999), which begins by running forward chaining on an unweighted version of the program.
Transformations on weighted logic programs / 55

of constit. The temporary item mentions all the capitalized variables in the expression.

Distributivity. A more important use appears when we combine folding with the distributive law. In the example above, the second rule’s body sums over the (underlined) free variables, J, Y, and Z. However, Y appears only in the temp item. We could therefore have summed over values of Y before multiplying by constit(Z,J,K), obtaining the following transformed program instead:

\[
\begin{align*}
temp2(X,Z,I,J) & += \text{rewrite}(X,Y,Z) \times \text{constit}(Y,I,J) \\
\text{constit}(X,I,K) & += \text{temp2}(X,Z,I,J) \times \text{constit}(Z,J,K).
\end{align*}
\]

This version of the transformation is permitted only because + distributes over \(*\). By “forgetting” Y as soon as possible, we have reduced the runtime of CKY from \(O(N^3 \cdot n^5)\) to \(O(N^3 \cdot n^2 + N^2 \cdot n^3)\).

Using the distributive law to improve runtime is a well-known technique. Aji and McEliece (2000) present what they call a “generalized distributive law,” which is equivalent to repeated application of the folding transformation. While their inspiration was the junction-tree algorithm for probabilistic inference in graphical models (discussed below), they demonstrate the approach to be useful on a broad class of weighted logic programs.

A categorial grammar view of folding. From a parsing viewpoint, notice that the item temp2(X,Z,I,J) can be regarded as a categorial grammar constituent: an incomplete X missing a subconstituent Z at its right (i.e., an X/Z) that spans the substring from I to J. This leads us to an interesting and apparently novel way to write the transformed program:

\[
\begin{align*}
\text{constit}(X,I,K)/\text{constit}(Z,J,K) & += \text{rewrite}(X,Y,Z) \times \text{constit}(Y,I,J) \\
\text{constit}(X,I,K) & += \text{constit}(X,I,K)/\text{constit}(Z,J,K) \times \text{constit}(Z,J,K).
\end{align*}
\]

Here A/B is syntactic sugar for slash(A,B). That is, / is used as an infix functor and does not denote division. However, it is meant to suggest division: as the second rule shows, A/B is an item which, if multiplied by B, yields a summand of A. In effect, the first rule above is derived from the original rule at the start of this section by dividing both sides by constit(Z,J,K). The second rule multiplies the missing factor constit(Z,J,K) back in, now that the first rule has summed over Y.

Notice that K appears free (and hence underlined) in the head of the first rule. The only slashed items that are actually provable by forward chaining are non-ground terms such as constit(s,0,K)/constit(vp,1,K). That is, they have the form constit(X,I,K)/constit(Z,J,K) where X,I,J are ground variables but K remains free. The way that K appears twice in the slashed item (i.e., internal

\footnote{All semiring-weighted programs enforce a similar distributive property. In particular, the trick can be applied equally well to common cases discussed in section 6.2.2: Viterbi parsing (\text{max} distributes over either \(*\) or \(\oplus\) and unweighted recognition (\(|\) distributes over \&).}
unification) indicates that the missing Z is always at the right of the X, while the fact that K remains a variable means that the shared right edge of the full X and missing Z are still unknown (and will remain unknown until the second rule fills in a particular Z). Thus, the first rule performs a computation once for all possible K—always the source of folding’s efficiency.

Our earlier program with temp2 could now be obtained by a further automatic transformation that replaces all \( \text{constit}(X,I,K)/\text{constit}(Z,J,K) \) having free K with the more compactly stored temp2(X,Z,I,J). The resulting rules are all range-restricted.

We emphasize that although our slashed items are inspired by categorial grammars, they can be used to describe folding in any weighted logic program. Section 6.5 will further exploit the analogy to obtain a novel “speculation” transformation.

**Further applications.** The folding transformation unifies various ideas that have been disparate in the natural language processing literature. Eisner and Satta (1999) speed up parsing with bilexical context-free grammars from \( O(n^3) \) to \( O(n^4) \), using precisely the above trick (see section 6.4 below). Huang et al. (2005) employ the same “hook trick” to improve the complexity of syntax-based MT with an n-gram language model.

Another parsing application is the common “dotted rule” trick (Earley, 1970). If one’s CFG contains ternary rules \( X \to Y_1 Y_2 Y_3 \), the naive CKY-like algorithm takes \( O(N^4 \cdot n^4) \) time:

\[
\text{constit}(X,I,L) \leftarrow (\text{rewrite}(X,Y_1,Y_2,Y_3) \cdot \text{constit}(Y_1,I,J)) \\
* \text{constit}(Y_2,J,K) \cdot \text{constit}(Y_3,K,L).
\]

Fortunately, folding allows one to sum first over \( Y_1 \) before summing separately over \( Y_2 \) and \( J \), and then over \( Y_3 \) and \( K \):

\[
\text{temp3}(X,Y_2,Y_3,I,J) \leftarrow \text{rewrite}(X,Y_1,Y_2,Y_3) \cdot \text{constit}(Y_1,I,J) \\
\text{temp4}(X,Y_3,I,K) \leftarrow \text{temp3}(X,Y_2,Y_3,I,J) \cdot \text{constit}(Y_2,J,K) \\
\text{constit}(X,I,L) \leftarrow \text{temp4}(X,Y_3,I,K) \cdot \text{constit}(Y_3,K,L).
\]

This restores \( O(n^3) \) runtime (more precisely, \( O(N^4 \cdot n^2 + N^3 \cdot n^3 + N^2 \cdot n^3) \))\(^{16}\) by reducing the number of nested loops. Even if we had declined to sum over \( Y_1 \) and \( Y_2 \) in the first two rules, then the summation over \( J \) would already have obtained \( O(n^3) \) runtime, in effect by binarizing the ternary rule. For example, temp4(X,Y1,Y2,Y3,I,K) would have corresponded to a partial constituent matching the dotted rule \( X \to Y_1 Y_2 \cdot Y_3 \). The additional summations over \( Y_1 \) and \( Y_2 \) result in a more efficient dotted rule that “forgets” the names of the nonterminals matched so far, \( X \to ? \cdot Y_3 \). This takes further advantage of distributivity by aggregating dotted-rule items (with +=) that will behave the same in subsequent computation.

\(^{16}\)For a dense grammar, which may have up to \( N^4 \) ternary rules. Tighter bounds on grammar size would yield tighter bounds on runtime.
The variable elimination algorithm for graphical models can be viewed as repeated folding. An undirected graphical model expresses a joint probability distribution over $P, Q$ by marginalizing (summing) over a product of clique potentials. In our notation,

$$\text{marginal}(P, Q) \leftarrow p_1(\ldots) \cdot p_2(\ldots) \cdot \ldots \cdot p_n(\ldots).$$

where a function such as $p_5(Q, X, Y)$ represents a clique potential over graph nodes corresponding to the random variables $Q, X, Y$. Assume without loss of generality that variable $X$ appears as an argument only to $p_{k+1}, p_{k+2}, \ldots, p_n$.

We may eliminate variable $X$ by transforming to

$$\text{temp}_5(\ldots) \leftarrow p_{k+1}(\ldots, X, \ldots) \cdot \ldots \cdot p_n(\ldots, X, \ldots).$$

$$\text{marginal}(P, Q) \leftarrow p_1(\ldots) \cdot \ldots \cdot p_k(\ldots) \cdot \text{temp}_5(\ldots).$$

Line 2 no longer mentions $X$ because line 1 has summed over it. To eliminate the remaining variables one at a time, the variable elimination algorithm applies this procedure repeatedly to the last line.\(^\text{17}\)

**Common subexpression elimination.** Folding can also be used multiple times to eliminate common subexpressions. Consider the following code, which is part of an inside algorithm for bilexical CKY parsing.\(^\text{18}\)

\[
\begin{align*}
\end{align*}
\]

Here $X:H$ is syntactic sugar for $\text{ntlex}(X,H)$, meaning a nonterminal $X$ whose head word is the lexical item $H$. The grammar uses two types of lexicalized binary productions (defined by $\text{rewrite}$ facts not shown here), which pass the head word to the left or right child, respectively.

We could fold together the last two factors of the first rule to obtain

\[
\begin{align*}
\text{temp}_6(Y:H,Z:H,I,K) & \leftarrow \text{constit}(Y:H,I,J) \cdot \text{constit}(Z:H,J,K). \\
\end{align*}
\]

We can reuse this definition of the temp rule to fold together the last two factors of line 3—which is the same subexpression, modulo variable renaming.

\(^\text{17}\)Determining the optimal elimination order is NP-complete. However, there are many heuristics in the literature (such as min-width) that could be used if automatic optimization of long rules is needed.

\(^\text{18}\)This algorithm is obviously an extension of the ordinary inside algorithm in section 6.2.2.

The other rules are

\[
\begin{align*}
\text{constit}(X:H,I,K) & \leftarrow \text{rewrite}(X,H) \cdot \text{word}(H,I,K). \\
\text{goal} & \leftarrow \text{constit}(s:H,0,N) \cdot \text{length}(N).
\end{align*}
\]
Given a new rule \( R \) in the form \( r \oplus F[s] \) (which will be used to replace a group of rules \( R_1, \ldots, R_n \) in \( P \)). Let \( S_1, \ldots, S_n \) be the complete list of rules in \( P \) whose heads unify with \( s \). Suppose that all rules in this list use \( \odot \) as their aggregation operator.

Now for each \( i \), when \( s \) is unified with the head of \( S_i \), the tuple \( (r, F_i, s, S_i) \) takes the form \( (r_i, F_i, s_i, s_i \odot = E_i) \). Suppose that for each \( i \), there is a distinct rule \( R_i \) in the program that is equal to \( r_i \odot = F_i[E_i] \), modulo renaming of its variables.

Then the folding transformation deletes the \( n \) rules \( R_1, \ldots, R_n \), and replaces them with the new rule \( R \), provided that

1. Any variable that occurs in any of the \( E_i \) which also occurs in either \( F_i \) or \( r_i \) must also occur in \( s_i \).
2. Either \( \odot = \) is simply \( = \) or else the distributive property \( [F[\mu] \odot F[\nu]] = [F[\mu \odot \nu]] \) holds for all assignments of terms to variables and all valuation functions \( [\cdot] \).

Before forming this 4-tuple, rename the variables in \( S_i \) so that they do not conflict with those in \( r, F, s \). Perform the desired unification within the 4-tuple by unifying it with the fixed term \( (R,F,S,S \odot = E) \), which contains two copies of \( S \).

This ensures that computing \( s_i \) by rule \( S_i \) does not sum over this variable, which would break the covariation of \( E_i \) with \( F \) or \( r \) as required by the original rule \( R_i \).

For instance, in the very first example of section 6.3, the \texttt{temp} item was defined using \( = \) and therefore performed no aggregation (see section 6.2.3). No distributivity was needed.

That is, all valuation functions over the space of ground terms, including dummy terms \( \mu \) and \( \nu \), when extended over expressions in the usual way.

\begin{figure}[h]
\centering
\includegraphics[width=\textwidth]{fig1.png}
\caption{The weighted folding transformation.}
\end{figure}

(\text{Below, for clarity, we explicitly and harmlessly swap the names of H2 and H within the temp rule.})

\begin{verbatim}
\end{verbatim}

Using the same \texttt{temp7} rule (modulo variable renaming) in both folding transformations, rather than introducing a new temporary item for each fold, gives us a constant-factor improvement in time and space.

\textbf{Formal definition of folding.} Our definition, shown in Figure 1, may seem surprisingly complicated. Its most common use is to replace a single rule \( r \oplus F[E] \) with \( r \oplus F[s] \) in the presence of a rule \( s \odot = E \). However, we have given a more general form that is best understood as precisely reversing the weighted unfolding transformation to be discussed in the next section (Figure 2). In unfolding, it is often useful for \( s \) to be defined by a group of rules
whose heads unify with $s$ (i.e., they may be more general or specific patterns than $s$). We define folding to allow the same flexibility.

In particular, this definition of folding allows an additional use of distributivity. Both the original item $r$ and the temp item $s$ may aggregate values not just within a single rule (summing over free variables in the body), but also across $n$ rules. In ordinary mathematical notation, we are performing a generalized version of the following substitution:

\[
\begin{align*}
\text{Before} & \quad r = \sum_{i=1}^{n} (f \ast E_i) \\
\text{After} & \quad r = f \ast s \\
\text{Before} & \quad s = \sum_{i=1}^{n} E_i \\
\text{After} & \quad s = \sum_{i=1}^{n} E_i
\end{align*}
\]

given the distributive property $\sum_i (f \ast E_i) = f \ast \sum_i E_i$. The common context in the original rules is the function “multiply by expression $f$,” so the temp item $s$ plays the role of $r/f$.

Figure 1 also generalizes beyond “multiply by $f$.” It allows an arbitrary common context $F$—a sort of function. In Figure 1 and throughout the paper, we use the notation $F[E]$ to denote the literal substitution of expression $E$ for all instances of $\mu$ in an expression $F$ over items, even if $X$ contains variables that appear in $E$ or elsewhere in the rule containing $F[E]$. We assume that $\mu$ is a distinguished symbol that does not appear elsewhere.

**Generalized distributivity.** Figure 1 states the required distributive property as generally as possible in terms of $F$. An interesting example is $[\log(p) + \log(q)] = [\log(p \ast q)]$, which says that $\log$ distributes over $\ast$ and changes it to $+$. This means that the definition $s(J) = e(J,K)$ may be used to replace $r += b(I,J) \ast \log(e(J,K))$ with $r += b(I,J) \ast \log(s(J))$. Here $n = 1$, $F = b(I,J) \ast \log(\mu)$, $E_1 = e(J,K)$, and $s = s(J)$.

By contrast, the definition $s += e(J)$ may not be used to replace $r += e(J) \ast e(J)$ with $r += s \ast s$, which would incorrectly replace a sum of squares with a square of sums. If we take $F$ to be $e(J) \ast \mu$ or $\mu \ast e(J)$, it is blocked by the first requirement in Figure 1 (variable occurrence). If we take $F$ to be $\mu \ast \mu$, it is blocked by the second requirement (distributivity).

**Introducing slashed definitions for folding.** Notice that Figure 1 requires the rules defining the temp item $s$ to be in the program already before folding occurs. If necessary, their presence may be arranged by a trivial definition introduction transformation that adds $\text{slash}(r, F) \Leftarrow E_i$ for each $i$, where $\text{slash}$ is a new functor not used elsewhere in $P$, and $\Leftarrow$ is chosen to ensure the required distributive property. We then take $s$ to be $\text{slash}(r, F)$ (or if one wants to use syntactic sugar, $r/F$).

Note that then $s_i$ will be $\text{slash}(r_i, F_i)$, which automatically satisfies the requirement in Figure 1 that certain variables that occur in $F_i$ or $r_i$ must also occur in $s_i$. This technique of introducing slashed items will reappear in section 6.5, where it forms a fundamental part of our speculation transformation.
Let $R$ be a rule in $\mathcal{P}$, given in the form $r \oplus F[s]$. Let $S_1, \ldots, S_n$ be the complete list of rules in $\mathcal{P}$ whose heads unify with $s$. Suppose that all rules in this list use $\odot$ as their aggregation operator.

Now for each $i$, when $s$ is unified with the head of $S_i$, the tuple $(r, F, s, S_i)^{23}$ takes the form $(r_i, F_i, s_i, s_i \odot E_i)$. Then the unfolding transformation deletes the rule $R$, replacing it with the new rules $r_i \oplus F_i[E_i]$ for $1 \leq i \leq n$. The transformation is allowed under the same two conditions as for the weighted folding transformation:

- Any variable that occurs in any of the $E_i$ which also occurs in either $F_i$ or $r_i$ must also occur in $s_i$.
- Either $\odot = \odot$ is simply $\oplus$, or else we have the distributive property $[F[\mu \odot F[v]] = [F[\mu \odot v]]$.

Before forming this tuple, rename the variables in $S_i$ so that they do not conflict with those in $r, F, s$.

**FIGURE 2** The weighted unfolding transformation.

If no operator $\odot$ can be found such that the distributive property will hold, and $n = 1$, then one can still use folding without the distributive property (as in the example that opened this section). In this case, introduce a rule temp($E_1) = E_1$, and take $s$ to be temp($E_1), which “memoizes” the value of expression $E_1$. Again, this satisfies the requirements of Figure 1.

### 6.4 Unfolding and Rule Elimination: Inlining Subroutines

**Unfolding.** The inverse of the folding transformation, called **unfolding** (Figure 2), replaces $s$ with its definition inside the rule body $r \oplus F[s]$. This definition may comprise several rules whose heads unify with $s$. If $s$ is regarded as a subroutine call, then unfolding is tantamount to inlining that call.

Recall that a folding transformation leaves the asymptotic runtime alone, or may improve it when combined with the distributive law. Hence unfolding makes the asymptotic runtime the same or worse. However, it may help the practical runtime by reducing overhead. (This is exactly the usual argument for inlining subroutine calls.)

An obvious example is program specialization. Consider the inside algorithm in section 6.2.2. If we take the second line,$\text{constit}(X, I, K) += \text{rewrite}(X, Y, Z) \ast \text{constit}(Y, I, J) \ast \text{constit}(Z, J, K)$.

and unfold $\text{rewrite}(X, Y, Z)$ inside it, then we replace the above rule with a set of rules, one for each binary production of the grammar (i.e., each rule whose head unifies with $\text{rewrite}(X, Y, Z)$):
constit(s,I,K) += 1 * constit(np,I,J) * constit(vp,J,K).
constit(np,I,K) += 0.5 * constit(det,I,J) * constit(n,J,K).

The resulting parser is specialized to the grammar, perhaps providing a constant-time speedup. It avoids having to look up the value of rewrite(s,np,vp) or rewrite(np,det,n) since these are now “hard-coded” into specialized inference rules. A compiled version can also “hard-code” pattern matching routines against specialized patterns such as constit(np,I,J); such pattern matches are used during forward or backward chaining to determine which rules to invoke.

Note that recursive calls can also be unfolded. For example, constit is recursively defined in terms of itself. If we unfold the constit(np,I,J) inside the first of the new rules above, we get

constit(s,I,K) += 1 * 0.5 * constit(det,I,L) * constit(n,L,J) * constit(vp,J,K).
constit(s,I,K) += 1 * 0.4 * word("Dumbo",I,J) * constit(vp,J,K).
constit(s,I,K) += 1 * 0.1 * word("flies",I,J) * constit(vp,J,K).

Unfolding is often a precursor to other transformations. For example, the pattern constit(vp,I,J) above can now be transformed to constit(vp,I,J) for more efficient storage. Furthermore, constant subexpressions like 1 * 0.5 can now be replaced in the source code by their values—a transformation that is known, not coincidentally, as constant folding. We will see another useful example of this unfold-refold pattern below, and yet another when we derive the (Eisner and Satta, 1999) algorithm in section 6.5.1.

Rule elimination. A practically useful transformation that is closely related to unfolding is what we call rule elimination (Figure 3). Rather than fully expanding one call to subroutine s, it removes one of the defining clauses of s and requires all of its callers to do the work formerly done by that clause.

This may change or eliminate the definition of s, so the transformation is not semantics-preserving. The advantage of changing the semantics is that if some s items become no longer provable, then it is no longer necessary to store them in the chart. Thus, rule elimination saves space. It also shares the advantages of unfolding—it can specialize a program, move unification to compile-time, eliminate intermediate steps, and serve as a precursor to other

---

26 A similar space savings—while preserving semantics—could be arranged simply by electing not to memoize these items, so that they are computed on demand rather than stored. Indeed, if we extend our formalism so that a program can specify what to memoize, then it is not hard to combine folding and unfolding to define a transformation that acts just like rule elimination (in that the callers are specialized) and yet preserves semantics. The basic idea is to fold together all of the other clauses that define s, then unfold all calls to s (which accomplishes the specialization), and finally declare that s (which is no longer called) should not be memoized (which accomplishes the space savings). However, we suppress the details as beyond the scope of this paper. Our main interest in rule elimination in this paper is to eliminate rules for temp items, whose semantics were introduced by a previous transformation and need not be preserved.
Let $S$ be a rule of $P$ to eliminate, with head $s$. Let $R_1, R_2, \ldots, R_n$ be a complete list of rules in $P$ whose bodies may depend on $s$. Suppose that each $R_i$ can be expressed in the form $r_i \oplus F_i[s_i]$, where $s_i$ is a term that unifies with $s$ and $F_i$ is an expression that is independent of $s$.

For each $i$, when $s_i$ is unified with the head of $S$, the tuple $(r_i, F_i, s_i, S)$ takes the form $(r_i', F_i', s_i', S')$. Then the rule elimination transform removes rule $S$ from $P$ and, for each $i$, adds the new rule $r_i' \equiv F_i'[E_i']$ (while also retaining $R_i$). The transformation is allowed under the same two conditions as for weighted folding and unfolding:

- Any variable that occurs in any of the $E_i'$ which also occurs in either $F_i'$ or $r_i'$ must also occur in $s_i'$.
- Either $\odot$ is simply $=$, or else we have the distributive property $[F_i'[\mu]] = [F_i[\mu \odot v]]$.

**Warning:** This transformation alters the semantics of ground terms that unify with $s$.

---

24 That is, the bodies of all other rules in $P$ must be independent of $s$. The notion of independence relies on the semantics of expressions, not on the particular program $P$. An expression $E$ is said to be independent of a term $s$ if for any two valuation functions on ground terms that differ only in the values assigned to groundings of $s$, the extensions of these valuation functions over expressions assign the same values to all groundings of $E$.

25 For example, suppose $s$ is $s(X,X)$. Then the rule $r(X) += s(X,Y) * t(Y)$ should be expressed as $r(X) += (\mu * t(Y))[s(X,Y)]$, while $r(X) \min = s(X,Y) * s(X,Y)$ should be expressed as $r(X) += (\mu * \mu)[s(X,Y)]$ and $r \min = 3[s(X,X)]$. However, $r(X) += s(X,Y) * s(Y,Z)$ cannot be expressed in the required form at all. We regard $\mu$ as a ground term in considering whether $F_i$ is independent of $s$.

---

FIGURE 3 The weighted rule elimination transformation.

To see the difference between rule elimination and unfolding, let us start with the same example as before, and selectively eliminate just the single binary production rewrite(np,det,n) = 0.5. In contrast to unfolding, this no longer replaces the original general rule

$$\text{constit}(X,I,K) += \text{rewrite}(X,Y,Z) * \text{constit}(Y,I,J) * \text{constit}(Z,J,K).$$

with a slew of specialized rules. Rather, it keeps the general rule but adds a specialization

$$\text{constit}(\text{np},I,K) += 0.5 * \text{constit}(\text{det},I,J) * \text{constit}(\text{n},J,K).$$

while deleting rewrite(np,det,n) = 0.5 so that it does not also feed into the general rule.

**A recursive example of rule elimination.** An interesting recursive example is shown below. The original program is in the first column. Eliminating
Transformations on weighted logic programs / 63

its second or third rule gives the program in the second or third column, respectively. Each of these changes damages the semantics of $s$, as warned, but preserves the value of $r$.\footnote{Since $r$ was defined to equal the (original) value of $s$, it provides a way to recover the original semantics of $s$. Compare the similar construction sketched in footnote 26.}

\begin{center}
\begin{tabular}{l|l}
$s += 1.$ & $s += 1.$ \\
$s += 0.5\cdot s.$ & $s += 0.5\cdot s.$ \\
$s += 0.5\cdot 1.$ & $s += 0.5\cdot 0.5\cdot s.$ \\
$r += s.$ & $r += s.$ \\
$r += 1.$ & $r += 0.5\cdot s.$ \\
\end{tabular}
\end{center}

\begin{tabular}{l|l}
$|s| = 2$, $|r| = 2$ & $|s| = 1$, $|r| = 2$ \\
\end{tabular}

\textbf{Unfolding or rule elimination followed by folding.}\footnote{Rule elimination can also be used after another transformation, such as speculation, to clean away unnecessary temp items. See footnote 42.} Recall the bilexical CKY parser given near the end of section 6.3. The first rule originally shown there has runtime $O(N^3 \cdot n^5)$, since there are $N$ possibilities for each of $X,Y,Z$ and $n$ possibilities for each of $I,J,K,H,H_2$. Suppose that instead of that slow rule, the original programmer had written the following folded version:

\begin{align*}
\end{align*}

\begin{align*}
\text{constit}(X:H,I,K) & += \text{temp}_8(X:H,Z:H_2,I,J) \cdot \text{constit}(Z:H_2,J,K).
\end{align*}

This partial program has asymptotic runtime $O\left(N^3 \cdot n^4 + N^2 \cdot n^5\right)$, and needs $O(N^2 \cdot n^4)$ space to store the items (rule heads) it derives.

By either unfolding the call to $temp_8$ or eliminating the $temp_8$ rule, we recover the first rule of the original program:

\begin{align*}
\end{align*}

This worsens the time complexity to $O\left(N^3 \cdot n^5\right)$. The payoff is that now we can refold this rule differently—either as follows (Eisner and Satta, 1999),

\begin{align*}
\end{align*}

\begin{align*}
\end{align*}

or alternatively as already shown in section 6.3 (whose $temp$ item had the additional advantage of being reusable). Either way, the asymptotic time complexity is now $O\left(N^3 \cdot n^4 + N^2 \cdot n^4\right)$—better than the original programmer’s version.

How about the asymptotic space complexity? If the first step used rule elimination rather than unfolding, then it actually eliminated storage for the $temp_8$ items, reducing the storage requirements from $O\left(N^2 \cdot n^4\right)$ to $O\left(N \cdot n^7\right)$.\footnote{Recall the bilexical CKY parser given near the end of section 6.3. The first rule originally shown there has runtime $O(N^3 \cdot n^5)$, since there are $N$ possibilities for each of $X,Y,Z$ and $n$ possibilities for each of $I,J,K,H,H_2$. Suppose that instead of that slow rule, the original programmer had written the following folded version:}

\begin{align*}
\end{align*}

\begin{align*}
\text{constit}(X:H,I,K) & += \text{temp}_8(X:H,Z:H_2,I,J) \cdot \text{constit}(Z:H_2,J,K).
\end{align*}
Regardless, the refolding step increased the space complexity back to the original programmer’s $O(N^2 \cdot n^4)$.

6.5 Speculation: Factoring Out Chains of Computation

In the most important contribution of this paper, we now generalize folding to handle unbounded sequences of rules, including cycles. This speculation transformation, which is novel as far as we know, is reminiscent of gap-passing in categorial grammar. It has many uses; we limit ourselves to two real-world examples.

6.5.1 Examples of the Speculation Transformation

Unary rule closure. Unary rule closure is a standard optimization on context-free grammars, including probabilistic ones (Stolcke, 1995). We derive it here as an instance of speculation. Suppose we begin with a version of the inside algorithm that allows unary nonterminal rules as well as the usual binary ones:

- $\text{constit}(X,I,K) += \text{rewrite}(X,W) \cdot \text{word}(W,I,K)$.
- $\text{constit}(X,I,K) += \text{rewrite}(X,Y) \cdot \text{constit}(Y,I,K)$.
- $\text{constit}(X,I,K) += \text{rewrite}(X,Y,Z) \cdot \text{constit}(Y,I,J) \cdot \text{constit}(Z,J,K)$.

Suppose that the grammar includes a unary rule cycle. For example, suppose that $\text{rewrite}($np,s$)$ and $\text{rewrite}($s,np$)$ are both provable. Then the values of $\text{constit}($np,I,K$)$ and $\text{constit}($s,I,K$)$ “feed into each other.” Under forward chaining, updating either one will cause the other to be updated; this process repeats until numerical convergence.\(^{29}\)

This computation is somewhat time-consuming—yet it is essentially the same for every $\text{constit}($np,I,K$)$ we may start with, regardless of the particular span I–K or the particular input sentence. We would prefer to do the computation only once, “offline.”

A difficulty is that the computation does incorporate the particular real value of the $\text{constit}($np,I,K$)$ that we start with. However, if we simply ignore this factor during our “offline” computation, we can multiply it in later when we get an actual $\text{constit}($np,I,K$)$. That is, we compute speculatively before the particular $\text{constit}($np,I,K$)$ and its value actually arrive.

In the transformed code below, $\text{temp}(X,X_0)$ represents the inside probability of building up a $\text{constit}(X,I_0,K_0)$ from a $\text{constit}(X_0,I_0,K_0)$ by a sequence of 0 or more unary rules. In other words, it is the total probability of all (possibly empty) unary-rewrite chains $X \rightarrow^* X_0$. While line 2 of the transformed

\(^{29}\)If we gave the Viterbi algorithm instead, with $\text{max}=\times$ in place of $+=\times$, then convergence would occur in finite time (at least for a PCFG, where all rewrite items have values in $[0,1]$). The same algorithm applies.
Transformations on weighted logic programs

The program still computes these items by numerical iteration, it only needs to compute them once for each $X,X_0$, since they are now independent of the particular span $I_0-K_0$ covered by these two constituents.\(^{30}\)

\[
\begin{align*}
temp(X_0,X_0) &\equiv 1. \\
temp(X,X_0) &\equiv rewrite(X,Y) \cdot temp(Y,X_0). \\
other(constit(X,I,K)) &\equiv rewrite(X,W) \cdot word(W,I,K). \\
other(constit(X,I,K)) &\equiv rewrite(X,Y,Z) \cdot constit(Y,I,J) \cdot constit(Z,J,K). \\
constit(X,I,K) &\equiv temp(X,X_0) \cdot other(constit(X_0,I,K)).
\end{align*}
\]

The $temp(s,np)$ item sums the probabilities of the infinitely many unary-rewrite chains $s \rightarrow^* np$, which build np up into s using only line 2 of the original program. Now, to get values like $constit(s,4,6)$ for a particular input sentence, we can simply sum finite products like $temp(s,np) \cdot other(constit(np,4,6))$, where $other(constit(np,4,6))$ sums up ways of building an $constit(np,4,6)$ other than by line 2 of the original program.\(^{31}\)

The semantics of this program, which can derive non-ground terms, fully defined by section 6.2.3. We omit further discussion of how it executes directly under forward chaining (section 6.2.4). However, note that the program could be transformed into a more conventional dynamic program by applying rule elimination to the first rule (the only one that is not range-restricted).\(^{32}\)

For efficiency, our formal transformation adds a “filter clause” to each of the temp rules:

\[
\begin{align*}
temp(X_0,X_0) &\equiv 1 \text{ needed only if } constit(X_0,I_0,K_0). \\
temp(X,X_0) &\equiv rewrite(X,Y) \cdot temp(Y,X_0) \text{ needed only if } constit(X_0,I_0,K_0). \\
other(constit(X,I,K)) &\equiv rewrite(X,W) \cdot word(W,I,K). \\
other(constit(X,I,K)) &\equiv rewrite(X,Y,Z) \cdot constit(Y,I,J) \cdot constit(Z,J,K). \\
constit(X,I,K) &\equiv temp(X,X_0) \cdot other(constit(X_0,I,K)). \\
constit(X_0,I,K) &\equiv 1 \cdot other(constit(X_0,I,K)).
\end{align*}
\]

\(^{30}\)We remark that the first $n$ steps of this iterative computation could be moved to compile time, by eliminating line 1 as discussed below, specializing line 2 to the grammar by unfolding $rewrite(X,Y)$, and then computing the series sums by alternately unfolding the $temp$ items and performing constant folding to consolidate each $temp$ item’s summands.

\(^{31}\)For example, it includes derivations where the np is built from a determiner and a noun, or directly from a word, but not where it is built directly from some s or another np. Excluding the last option prevents double-counting of derivations.

\(^{32}\)Here is the result, which alters the semantics of the slashed $temp$ item to ignore derivations of length 0:

\[
\begin{align*}
temp(X,X_0) &\equiv rewrite(X,Y) \cdot temp(Y,X_0). \\
temp(X,X_0) &\equiv rewrite(X,X_0) \cdot 1. \\
other(constit(X,I,K)) &\equiv rewrite(X,W) \cdot word(W,I,K). \\
other(constit(X,I,K)) &\equiv rewrite(X,Y,Z) \cdot constit(Y,I,J) \cdot constit(Z,J,K). \\
constit(X,I,K) &\equiv temp(X,X_0) \cdot other(constit(X_0,I,K)). \\
constit(X_0,I,K) &\equiv 1 \cdot other(constit(X_0,I,K)).
\end{align*}
\]
are proved or updated only once some \(\text{constit}(X_0,I_0,K_0)\) constituent has been built.\(^{33}\) At that time, all the \(\text{temp}(X,X_0)\) values for this \(X_0\) will be computed once and for all, since there is now something for them to combine with. Usefully, these values will then remain static while the grammar does, even if the sentence changes.

Adopting the categorial view we introduced in section 6.3, we can regard \(\text{temp}(s,np)\) as merely an abbreviation for the \textit{non-ground} slashed item \(\text{constit}(s,I_0,K_0)/\text{constit}(np,I_0,K_0)\): the cost of building up a \(\text{constit}(s,I_0,K_0)\) if we already had a \(\text{constit}(np,I_0,K_0)\). This cost is independent of \(I_0\) and \(K_0\), which is why we only need to compute a single item to hold it, albeit one that contains variables.

As we will see, the slashed notation is not merely expository. Our formal speculation transformation will actually produce a program with slashed terms, essentially as follows:

\[
\begin{align*}
\text{constit}(X_0,I_0,K_0)/\text{constit}(X_0,I_0,K_0) &\quad += 1. \\
\text{constit}(X,I,K)/\text{constit}(X_0,I_0,K_0) &\quad += \text{rewrite}(X,Y) * \text{constit}(Y,I,K)/\text{constit}(X_0,I_0,K_0). \\
\text{other}(\text{constit}(X,I,K)) &\quad += \text{rewrite}(X,W) * \text{word}(W,I,K). \\
\text{other}(\text{constit}(X,I,K)) &\quad += \text{rewrite}(X,Y,Z) * \text{constit}(Y,I,J) * \text{constit}(Z,J,K). \\
\text{constit}(X,I,K) &\quad += (\text{constit}(X,I,K)/\text{constit}(X_0,I_0,K_0)) * \text{other}(\text{constit}(X_0,I_0,K_0)).
\end{align*}
\]

\textbf{A variable-free example.} To understand better how the slash and other mechanisms work, consider this artificial variable-free program, illustrated by the hypergraph in Figure 4:

\[
\begin{align*}
a &\quad += b * c. \\
b &\quad += r. \\
c &\quad += f * c. \\
c &\quad += d * e * x. \\
c &\quad += g. \\
x &\quad += \ldots
\end{align*}
\]

The values of \(a\) and \(c\) depend on \(x\). We elect to create speculative versions of the first, third, and fourth rules. The resulting program is drawn in Figure 5. It includes rules to compute slashed versions of \(a, c\) and \(x\) itself that are “missing an \(x\)”:

\[
\begin{align*}
a/x &\quad += b * c/x. \\
c/x &\quad += f * c/x. \\
c/x &\quad += d * e * x/x. \\
x/x &\quad += 1.
\end{align*}
\]

\(^{33}\text{In this example, the filter clause on the second rule is redundant. Runtime analysis or static analysis could determine that it has no actual filtering effect, allowing us to drop it.}\)
FIGURE 4 A simple variable-free program before applying the speculation transformation.

FIGURE 5 The program of Figure 4 after applying the speculation transformation. The $x/x$ rule and various $\text{other(...)}$ rules have been eliminated for simplicity.
It also reconstitutes full-fledged versions of \( a, c, \) and \( x \). Each is defined by a sum that is split into two cases: summands that were built from an \( x \) using a sequence of \( \geq 0 \) of the selected rules, and “other” summands that were not. (Notice that the first rule is not \( a = a/x \ast x \); this is because \( x \) might in general include derivations that are built from another \( x \) (though not in this example), and this would lead to double-counting. By using \( a = a/x \ast \text{other}(x) \), we split each derivation of \( a \) uniquely into a maximal sequence of selected rules, applied to a minimal instance of \( x \).)

\[
\begin{align*}
a & \leftarrow a/x \ast \text{other}(x). \\
c & \leftarrow c/x \ast \text{other}(x). \\
x & \leftarrow x/x \ast \text{other}(x). \\
a & \leftarrow \text{other}(a). \\
c & \leftarrow \text{other}(c).
\end{align*}
\]

Finally, the program must define the “other” summands:

\[
\begin{align*}
\text{other}(a) & \leftarrow b \ast \text{other}(c). \\
\text{other}(c) & \leftarrow f \ast \text{other}(c). \\
\text{other}(x) & \leftarrow g.
\end{align*}
\]

In Figure 5, this program has been further simplified by eliminating the rules for \( x/x \) and \( \text{other}(x) \).

**Split bilexical grammars.** For our next example, consider a “split” bilexical CFG, in which a head word must combine with all of its right children before any of its left children. The naive algorithm for bilexical context-free parsing is \( O(n^5) \). In the split case, we will show how to derive the \( O(n^4) \) and \( O(n^3) \) algorithms of Eisner and Satta (1999).

The “inside algorithm” below builds up \( \text{rconstit} \) items by successively adding 0 or more child constituents to the right of a word, then builds up \( \text{constit} \) items by adding 0 or more child constituents to the left of this \( \text{rconstit} \). As before, \( X:H \) represents a nonterminal \( X \) whose head word is \( H \).

\[
\begin{align*}
\text{rconstit}(X:H,I,K) & \leftarrow \text{rewrite}(X,H) \ast \text{word}(H,I,K). \% 0 \text{ right children so far} \cr
\text{rconstit}(X:H,I,K) & \leftarrow \text{rewrite}(X:H,Y:H,Z:H2) \ast \text{rconstit}(Y:H,I,J,K) \ast \text{constit}(Z:H2,J,K). \% \text{add right child} \cr
\text{constit}(X:H,I,K) & \leftarrow \text{rconstit}(X:H,I,K). \% \text{0 left children so far} \cr
\text{constit}(X:H,I,K) & \leftarrow \text{rewrite}(X:H,Y:H2,Z:H) \ast \text{constit}(Y:H2,I,J) \ast \text{constit}(Z:H,J,K). \% \text{add left child} \cr
\text{goal} & \leftarrow \text{constit}(s:H,0,N) \ast \text{length}(N).
\end{align*}
\]

\[^{34}\text{We deal here with context-free grammars, rather than the head-automaton grammars of Eisner and Satta. In particular, our complete constituents carry nonterminal categories and not just head words. Note that the algorithm is correct only for a “split” grammar (formally, one that does not contain two rules of the form } \text{rewrite}(X:H1,Y:H2,Z:H1) \text{ and } \text{rewrite}(V:H1,X:H1,W:H3), \text{ since otherwise the grammar would license trees that cannot be constructed by the algorithm.}\]
This obvious algorithm has runtime $O(N^3 \cdot n^5)$ (dominated by line 4). We now speed it up by exploiting the conditional independence of left children from right children. To build up a $\text{constit}$ whose head word starts at $I$, we will no longer start with a particular, existing $\text{rconstit}$ from $I$ to $K$ (line 3) and then add left children (line 4). Rather, we transform the program so that it abstracts away the choice of starting $\text{rconstit}$. It can then build up the $\text{constit}$ item speculatively, adding left children without having committed to any particular $\text{rconstit}$. As this work is independent of the $\text{rconstit}$, the items derived during it do not have to specify any value for $K$. Thus, work is shared across all values of $K$, improving the asymptotic complexity. Only after finishing this speculative computation does the program fill in each of the various $\text{rconstit}$ items that could have been chosen at the start. To accomplish this transformation, replace lines 3–4 with

\begin{align*}
\text{lconstit}(X_0:H_0,X_0,I_0,J_0) &\leftarrow 1. \\
\text{lconstit}(X:H_0,X_0,I,J_0) &\leftarrow \text{rewrite}(X:H_0,Y:H_2,Z:H) \\
&\quad \ast \text{constit}(Y:H_2,I,J_0) \ast \text{lconstit}(Z:H_0,X_0,J_0). \\
\text{constit}(X:H_0,I,K_0) &\leftarrow \text{lconstit}(X:H_0,X_0,I,J_0) \ast \text{rconstit}(X_0:H_0,J_0,K_0).
\end{align*}

The new temp item $\text{lconstit}(X:H_0,X_0,I_0,J_0)$ represents the left half of a constituent, stretching from $I$ to $J_0$. We can regard it again in categorial terms: as the last line suggests, it is just a more compact notation for a $\text{constit}$ missing its $\text{rconstit}$ right half from $J_0$ to some $K_0$. This can be written more perspicuously as $\text{constit}(X:H_0,I,K_0)/\text{rconstit}(X_0:H_0,J_0,K_0)$, where $K_0$ is always a free variable, so that $\text{lconstit}$ need not specify any particular value for $K_0$.

The first $\text{lconstit}$ rule introduces an empty left half. This is extended with its left children by recursing through the second $\text{lconstit}$ rule, allowing $X$ and $I$ to diverge from $X_0$ and $J_0$ respectively. Finally, the last rule fills in the missing right half $\text{rconstit}$.

Again, our speculation transformation will actually produce the slashed notation as its output. Specifically, it will replace lines 3–4 of the original untransformed program with the following.\(^{35}\)

\begin{align*}
\text{rconstit}(X_0:H_0,J_0,K_0)/\text{rconstit}(X_0:H_0,J_0,K_0) &\leftarrow 1 \\
\text{needed\_only\_if} \quad \text{rconstit}(X_0:H_0,J_0,K_0). \\
\text{constit}(X:H,J,K)/\text{rconstit}(X_0:H_0,J_0,K_0) &\leftarrow \text{rconstit}(X:H,J,K)/\text{rconstit}(X_0:H_0,J_0,K_0) \\
&\quad \text{needed\_only\_if} \quad \text{rconstit}(X_0:H_0,J_0,K_0).
\end{align*}

\(^{35}\)In fact our transformation in Figure 6 will produce something a bit more complicated. The version shown has been simplified by using rule elimination (section 6.4) to trim away all other (…) items, which do not play a significant role in this example. That is because the only slashed items are $\text{constit}/\text{rconstit}$, and there is no other way to build a $\text{constit}$ except from an $\text{rconstit}$.
\begin{align*}
\text{constit}(X:H,I,K)/r\text{constit}(X0:H0,J0,K0) \\
& \quad \quad \rightarrow \text{rewrite}(X:H,Y:H2,Z:H) \ast \text{constit}(Y:H2,I,J) \\
& \quad \quad \quad \ast \text{constit}(Z:H,J,K)/r\text{constit}(X0:H0,J0,K0) \\
& \quad \quad \text{needed, only if } r\text{constit}(X0:H0,J0,K0).
\end{align*}

\begin{align*}
\text{constit}(X:H,I,K) & \quad \rightarrow \text{constit}(X:H,I,K)/r\text{constit}(X0:H0,J0,K0) \\
& \quad \quad \ast r\text{constit}(X0:H0,J0,K0).
\end{align*}

The first line introduces a slashed item. The next two lines are the result of slashing \(r\text{constit}(X0:H0,J0,K0)\) out of the original lines 3–4; note that \(X0, H0, J0,\) and \(K0\) appeared nowhere in the original program. The final line reconstitutes the \text{constit} item defined by the original program, so that the transformed program preserves the original program’s semantics.

By inspecting this program, one can see that the only provable items of the form \(\text{constit}(X:H,I,K)/r\text{constit}(X0:H0,J0,K0)\) actually have \(H=H0, K=K0,\) and \(K0\) a free variable.\(^{36}\) These conditions are true for the slashed item that is introduced in the first line, and they are then preserved by every rule that derives a new slashed item. This is why in our earlier presentation of this code, we were able to abbreviate such a slashed item by \(l\text{constit}(X:H0,X0,I,J0)\), which uses only 5 variables rather than 8. Discovering such abbreviations by static analysis is itself a transformation that we do not investigate in this paper.

**Filter clauses can improve asymptotic runtime.** The special filter clause \( \text{needed, only if } r\text{constit}(X0:H0,J0,K0) \) is added solely for efficiency, as always. It says that it is not necessary to build a left half that might be useless (i.e., purely speculatively), but only when there is at least one right half for it to combine with.

In this example, the filter clause is subtly responsible for a voiding an extra factor of \(V\) in the runtime, where \(V \gg n\) is the size of the vocabulary. For simplicity, let us return to the unslashed notation:

\begin{align*}
l\text{constit}(X:H0,X0,I,J0) & \quad \rightarrow \text{rewrite}(X:H0,Y:H2,Z:H) \\
& \quad \quad \ast \text{constit}(Y:H2,I,J) \ast l\text{constit}(Z:H0,X0,J,J0). \\
& \quad \quad \text{needed, only if } r\text{constit}(X0:H0,J0,K0).
\end{align*}

The intent is to build only left halves \(l\text{constit}(H:H0,X0,I,J0)\) whose head word \(H0\) will actually be found starting at the right edge \(J0\). However, without the filter, the above rule could be far more speculative, combining a finished left child such as \(\text{constit}(np:dumbo,4,5)\) with a rewrite rule such as \(\text{rewrite}(s:flies,np:dumbo,vp:flies)\) and the non-ground item \(l\text{constit}(X0:H0,X0,J0,J0)\) (defined elsewhere with value 1) to obtain \(l\text{constit}(s:flies,vp,4,5)\)—regardless of whether \(flies\) actually starts at position 5 or even appears in the input sentence at all! This would lead to a proliferation of

\(^{36}\)By contrast, we already noted that \(X\) and \(I\) could diverge from \(X0\) and \(J0\) respectively, in this particular program.
$O(V) \text{ lconstit items with speculative head words such as flies that might start at position } 5$. The filter clause prevents this by “looking ahead” to see whether any items of the form rconstit(vp:flies,5,K0) have actually been proved.

As a result, the runtime is now $O(n^4)$ (as compared to $O(n^5)$ for the un-transformed program).\footnote{We could also have achieved $O(n^4)$ simply by folding the original program as discussed in section 6.3. However, that would not have allowed the further reduction to $O(n^3)$ discussed below.} This is so because the rule above may be grounded in $O(n^4)$ ways reflecting different bindings of I, J, J0, and word H2, where H2 may in practice be any of the words in the span I–J. Although the rule also mentions H0, the filter clause has ensured that H0’s binding is completely determined by J0’s.

As a bonus, we can now apply the unfold-refold pattern to obtain the $O(n^3)$ algorithm of Eisner and Satta (1999). Starting with our transformed program, unfold constit in the body of each rule where it appears,\footnote{Including if desired the goal rule, not discussed here. The old constit rule is then useless, except perhaps to the user, and may be trimmed away if desired by rule elimination.} giving

\begin{align*}
\text{rconstit} &\leftarrow \text{rconstit} \ast \text{rewrite} \ast (\text{lconstit} \ast \text{rconstit}). \\
\text{lconstit} &\leftarrow (\text{lconstit} \ast \text{rconstit}') \ast \text{rewrite} \ast \text{lconstit}.
\end{align*}

where the ‘ symbol marks the halves of the unfolded constit, and the three adjacent half-constituents are written in left-to-right order. Now re-parenthesize these rules as

\begin{align*}
\text{rconstit} &\leftarrow (\text{rconstit} \ast (\text{rewrite} \ast \text{lconstit}')) \ast \text{rconstit}'. \\
\text{lconstit} &\leftarrow (\text{lconstit}' \ast (\text{rconstit}' \ast \text{rewrite}) \ast \text{lconstit}).
\end{align*}

and fold out each parenthesized subexpression, using distributivity to sum over its free variables. The items introduced when folding the large subexpressions correspond, respectively, to Eisner and Satta’s “right trapezoid” and “left trapezoid” items. The speedup arises because there are $O(n)$ fewer possible trapezoids than constits: a constit had to specify a head word that could be any of the words covered by the constit, but a trapezoid’s head word is always the word at its left or right edge.

### 6.5.2 Semantics and Operation of Filter Clauses

Our approach to filtering is novel. Our \texttt{needed-only} clauses may be regarded as “relaxed” versions of side conditions (Goodman, 1999). In the denotational semantics (section 6.2.3), they relax the restrictions on the valuation function, allowing more possible valuations for the transformed program. (In the case of speculation, these valuations may disagree on the new slashed items, but all of them preserve the semantics of the original program.)

Specifically, when constructing $P(r)$ to determine whether a ground item $r$ is provable and what its value is, we may \textit{optionally} omit the summand corre-
responding to a grounded rule \( r \oplus = E \) if this rule has an attached filter clause needed\_only\_if \( C \) such that no consistent grounding of \( C \) has been proved.\(^\text{39}\)

How does this help operationally, in the forward chaining algorithm? When a rule triggers an update to a ground or non-ground item, but carries a (partly instantiated) filter clause that does not unify with any proved item, then the update has infinitely low priority and need not be propagated further by forward chaining. The update must still be carried out if the filter clause is proved later.

The optionality of the filter is crucial for two reasons. First, if a filter becomes false, the forward-chaining algorithm is not required to retract updates that were performed when the filter was true. In the examples of section 6.5.1 above or section 6.6.3 below, the filter clauses ensure that entries are filled into the unary-rule-closure and left-corner tables only as needed. Once this work has been done, however, these entries are allowed to persist even when they are no longer needed, e.g., once the facts describing the input sentence are retracted. This means that we can reuse them for a future sentence rather than re-deriving them every time they are needed.

Second, the forward-chaining algorithm is not required to bind variables in the rule when it checks for consistent groundings of the filter clause. Consider this rule from earlier:

\[
\text{constit}(X:H,1,K)/r\text{constit}(X0:H0,J0,K0) \\
\quad \quad \rightarrow \text{rewrite}(X:H,Y:H2,Z:H) \ast \text{constit}(Y:H2,I,J) \\
\quad \quad \ast \text{constit}(Z:H,J0,K)/r\text{constit}(X0:H0,J0,K0) \\
\quad \text{needed\_only\_if} \ r\text{constit}(X0:H0,J0,K0).
\]

Recall that the only \text{constit}/r\text{constit} items that are actually derived are non-ground items in which \( K=K0 \) and is free, such as \text{constit}(s:flies,4,K0)/r\text{constit}(vp:flies,5,K0). Such a non-ground item actually represents an infinite collection of possible ground items that specialize it. The semantics of needed\_only\_if, which are defined over ground terms, say that we only need to derive a subset of this collection: rather than proving the non-ground item above, we are welcome to prove only the “needed” ground instantiations, with specific values of \( K0 \) such that \text{rconstit}(vp:flies,5,K0) has been proved. However, in general, this would prove \( O(n) \) ground items rather than a single non-ground item. It would destroy the whole point of speculation, which is to achieve a speedup by leaving \( K0 \) free until specific \text{rconstit} items are multiplied back in at the end. Thus, the forward-chaining algorithm is better off exploiting the optionality of filtering and proving the non-ground version—thus

\(^{39}\)The “consistent” groundings are those in which variables of \( C \) that are shared with \( r \) or \( E \) are instantiated accordingly. In the speculation transformation, all variables of \( C \) are in fact shared with \( r \) and \( E \). If they were not, \( C \) could have many consistent groundings, but we would still aggregate only one copy of \( E \), just as if the filter clause were absent, not one per copy per consistent grounding.
proving more than is strictly needed—as long as at least one of its groundings is needed (i.e., as long as some item that unifies with \texttt{rconstit(vp:flies,4,K0)} has already been proved).

For a simpler example, consider

\begin{verbatim}
\texttt{rconstit(X0:H0,J0,K0)/rconstit(X0:H0,J0,K0) += 1}
\end{verbatim}

\texttt{needed only if rconstit(X0:H0,J0,K0)}. 

A reasonable implementation of forward chaining will prove the non-ground item \texttt{rconstit(X0:H0,J0,K0)/rconstit(X0:H0,J0,K0)}—just as if the filter clause were not present—provided that \textit{some} grounding of \texttt{rconstit(X0:H0,J0,K0)} has been proved. It is not required to derive a separate grounding of the slashed item for each grounding of \texttt{rconstit(X0:H0,J0,K0)}; this would also be correct but would be less efficient.

### 6.5.3 Formalizing Speculation for Semiring-Weighted Program Fragments

To formalize the speculation transformation, we begin with a useful common case that suffices to handle our previous examples. This definition (Figure 6) allows us to speculate over a set of rules that manipulate values in a semiring of weights \( W \). Such rules must all use the same aggregation operator, which we call \( \oplus \), with identity element \( 0 \). Furthermore, each rule body must be a simple product of one or more items, using an associative binary operator \( \otimes \) that distributes over \( \oplus \) and has an identity element \( 1 \). This version of the transformation is only able to slash the final term in such a product; however, if \( \otimes \) is commutative, then the terms in a product can be reordered to make any term the final term.

Our previous examples of speculation can be handled by taking the semiring \((W, \oplus, \otimes, 0, 1)\) to be \((\mathbb{R}, +, \cdot, 0, 1)\). Moreover, any unweighted program can be handled by taking \((W, \oplus, \otimes, 0, 1)\) to be \((\{F, T\}, |, \&, F, T)\).

Intuitively, \texttt{other(A)} in Figure 6 accumulates ways of building \( A \) other than groundings of \( F_i \otimes F_j \otimes \cdots \otimes F_{j_0} \otimes x \) for \( j > 0 \). Meanwhile, \texttt{slash(A,x)} accumulates ways of building \( A \) by grounding products of the form \( F_i \otimes F_j \otimes \cdots \otimes F_{j_0} \) for \( j \geq 0 \). To “fill in the gap” and recover the final value of \( A \) (as is required to preserve semantics), we multiply \texttt{slash(A,x)} only by \texttt{other(x)} (rather than by \texttt{x}), in order to prevent double-counting (which is analogous to spurious ambiguity in a categorial grammar).

To apply our formal transformation in the unary-rule elimination example, take \( x=\texttt{constit(X0,0,0,0)} \). As required, \( X0,0,0,0 \) do not appear in the original program. Take \( R_1 \) to be the “unary” \texttt{constit} rule of the original program where \( t_i \) is the last item in the body of \( R_i \). Here \( k = 0 \).

To apply our formal specification in the artificial variable-free example of Figures 4–5, take \( x = x \), \( R_1 = a += b \cdot c \), \( R_2 = c += f \cdot c \), and \( R_3 = c += (d \cdot e) \cdot x \). Since, among the \( t_i \), only \( t_3 \) unifies with \( x \), we have \( k = 2 \).
Given a semiring \((\mathcal{W}, \circ, \o, \text{true})\).
Given a term \(x\) to slash out, where any variables in \(x\) do not occur anywhere in the program \(P\). Given distinct rules \(R_1, \ldots, R_n\) in \(P\) from which to simultaneously slash \(x\) out, where each \(R_i\) has the form \(r_i \circ F_i \circ t_i\) for some expression \(F_i\) (which may be \(\text{true}\)) and some item \(t_i\).
Let \(k\) be the index\(^{40}\) such that \(0 \leq k \leq n\) and
- For \(i \leq k\), \(t_i\) does not unify with \(x\).
- For \(i > k\), \(t_i\) unifies with \(x\); moreover, their unification equals \(t_i\).

Then the speculation transformation constructs the following new program. Recall that \(\oplus\) denotes the aggregation operator for \(r\) (which may or may not be \(\circ\)). Let slash, other, and matches\(_x\) be new functors that do not already appear in \(P\).

- \(\text{slash}(x,x) \circ \oplus = \text{true}\) needed\(_x\).
- \((\forall 1 \leq i \leq n) \text{slash}(r_i, x) \circ \oplus = F_i \circ \text{slash}(t_i, x)\) needed\(_x\).
- \((\forall 1 \leq i \leq k) \text{other}(r_i) \circ \oplus = F_i \circ \text{other}(t_i)\).
- \((\forall \text{ rules } p \circ \oplus = E \text{ not among the } R_i) \text{ other}(p) \circ \oplus = E\).
- \(\text{matches\(_x\)}(x) \models \text{true}\). \(\text{matches\(_x\)}(A) \models \text{false}\).
- \(A \circ \oplus = \text{other}(A)\) if not matches\(_x\)(A).\(^{43}\)
- \(A \circ \oplus = \text{slash}(A,x) \circ \text{other}(x)\).

\(^{40}\)If necessary, the program can be pre-processed so that such an index exists. Any rule can be split into three more specialized rules: an \(i \leq k\) rule, an \(i > k\) rule, and a rule not among the \(R_i\). Some of these rules may require boolean side conditions to restrict their applicability.

\(^{41}\)That is, \(i\) is “more specific” than \(x\): it matches a non-empty subset of the ground terms that \(x\) does.

\(^{42}\)It is often worth following the speculation transformation with a rule elimination transformation (section 6.4) that removes some of these \(\text{other}\) items. In particular, if \(p\) does not unify with \(x\) or any of the \(t_i\), then the only rule that uses \(\text{other}(p)\) is \(A \circ \oplus = \text{other}(A)\). In this case, eliminating the old rule \(\text{other}(p) \circ \oplus = E\) simply restores the original rule \(p \circ \oplus = E\).

\(^{43}\)Note that \(A\) ranges over all ground terms. (Except those that unify with \(x\), which are covered by the next rule.) The aggregation into a particular ground term \(A\) must be handled using the appropriate aggregation operator for that ground term, here denoted \(\oplus = \circ \text{other}(\oplus)\). In the example programs, this awkward notation was avoided by splitting this rule into several rules, which handle various disjoint classes of items \(A\) that have different aggregation operators.

**FIGURE 6** The semiring-weighted speculation transformation.

To apply our formal transformation in the split bilexical grammar example, take \(x=\text{constit}(X0:H0,J0,K0)\), the \(R_i\) to be the two rules defining \(\text{constit}\), each \(t_i\) to be the last item in the body of \(R_i\), and \(k = 1\).

**Folding as a special case of speculation.** As was mentioned earlier, the folding transformation is a special case of the speculation transformation, in which application is restricted to rules with a single ground term at their head,
and the item to be slashed out must appear literally in each affected rule. For ease of presentation, however, the formulations above are not quite parallel. In folding, we adopt the convention that a common function $F$ is being “slashed out” of a set of rules, and the different items to which that function applies are aggregated first into a new intermediate item. In speculation, we take the opposite view, where there is a common item to be slashed out present as the argument to different functions, so that the functions get aggregated into a new lambda term. We chose the former presentation for folding to avoid the needless complication of using the lambda calculus, but we needed the flexibility of the latter for a fully general version of speculation. In the case of ordinary semiring-weighted programs, this distinction is trivial; when slashing out item $a$ from a rule like $f \leftarrow a \cdot b$, we can equally easily say that we are slashing out the function “multiply by $a$” from its argument $b$ or that we are slashing out the item $a$ from inside the function “multiply by $b$”. However, in general, being able to leave behind functions allows us to construct intermediate terms which don’t carry a numerical value; for example, we could choose to slash out the $a$ item from a rule like $f \leftarrow \log(a)$ and propagate just the function $\text{slash}(f, a) \leftarrow \lambda x . \log(x)$.

6.5.4 Formalizing Speculation for Arbitrary Weighted Logic Programs

The speculation transformation becomes much more complicated when it is not restricted to semiring-weighted programs. In general, the value of a slashed item is a function, just like the semantics of a slashed constituent in categorial grammar. Functions are aggregated pointwise: that is, we define $(\lambda z . f(z)) \odot (\lambda z . g(z)) = \lambda z . (f(z) \odot g(z))$.

As in categorial grammar semantics, gaps are introduced with the identity function, passed with function composition, and eliminated with function application.

In the commutative semiring-weighted programs discussed above, all functions had the form “multiply by $w$” for some weight $w$. We were able to avoid the lambda-calculus there by representing such a function simply as $w$, and by using $1$ for the identity function, semiring multiplication $\odot$ for both composition and application, and semiring addition $\oplus$ for pointwise addition.

We defer the details of the formal transformation to a later paper. It is significantly more complicated than Figure 6 because we can no longer rely on the mathematical properties of the semiring. As in folding and unfolding (Figures 1–2), we must demand a kind of distributive-law property to ensure that the semantics will be preserved (recall the log example from section 6.3). This property is harder to express for speculation, which is like folding through unbounded sequences of rules, including cycles.

Consider the semiring-weighted program in Figures 4–5. The original program only used the item $x$ early in the computation, multiplying it by $d \cdot e$. 
The transformed program had to reconstitute \( a \) from \( a/x \) and \( x \) (and \( c \) from \( c/x \) and \( x \)). This meant multiplying \( x \) in later, only after the original \( d \ast e \) had passed through several levels of \( \ast \) and \( += \) in the rule \( a += b \ast c \) and the cyclic rule \( c += f \ast c \).

In general, we want to be sure that delaying the introduction of \( x \) until after several intermediate functions and aggregations does not change the value of the result. Hence, a version of the distributive property must be enforced at each intermediate rule affecting the slashed items.

Furthermore, if the slashed-out item \( x \) contains variables, then introducing it will aggregate over those variables. For example, the rule \( a(B,C) += (\text{a}(B,C)/x(B0,C0,D0))(x(B0,C0,D0)) \) not only applies a function to an argument, but also aggregates over \( B0 \), \( C0 \), and \( D0 \). In the original version of the program, these aggregations might have been performed with various operators. When \( a(B,C) \) is reconstituted in the transformed version, we must ensure that the same sequence of aggregations is observed. In order to do this correctly, it is necessary to keep track of the association between the variables being aggregated in the original program and the variables in the slashed item, so that we can ensure that the same aggregations are performed.

### 6.6 Magic Templates: Filtering Useless Items

The bottom-up “forward-chaining” execution strategy of section 6.2.4 will (if it converges) compute values for all provable items. Typically, however, the user is primarily interested in certain items (often just goal) and perhaps the other items involved in deriving them.

In parsing, for example, the user may not care about building all legal constituents—only those that participate in a full parse. Similarly, in a program produced by speculation, the user does not care about building all possible slashed items \( r/x \)—only those that can ultimately combine with some actual \( x \) to reconstitute an item \( r \) of the original program.

Can we prevent forward chaining from building at least some of the “useless” items? In the speculation transformation (section 6.5 and Figure 6), we accomplished this by filtering our derivations with \text{needed\_only\_if} \text{clauses.}

We now give a transformation that explains and generalizes this strategy. It prevents the generation of some “useless” items by automatically adding \text{needed\_only\_if} filter clauses to an existing program. A version of this \text{magic templates} transformation was originally presented in a well-known paper by Ramakrishnan (1991), generalizing an earlier transformation called magic sets.

#### 6.6.1 An Overview of the Transformation

Since this transformation makes some terms unprovable, it cannot be semantics-preserving. We will say that a ground term \( a \) is \text{charmed} if we
should preserve its semantics. In other words, the semantic valuation functions of the transformed program and the original program will agree on at least the charmed terms. The program will determine at runtime which terms are charmed: a ground term \( a \) is considered charmed if the term \( \text{magic}(a) \) is provable (inevitably with value true).

The user should explicitly charm the ground terms of interest to him or her by writing rules such as

\[
\begin{align*}
\text{magic}(\text{goal}) & \rightarrow \text{true}. \\
\text{magic}(\text{rewrite}(X,Y,Z)) & \rightarrow \text{true}. \\
\text{magic}(\text{word}(W,I,J)) & \rightarrow \text{true}. \\
\text{magic}(\text{length}(N)) & \rightarrow \text{true}.
\end{align*}
\]

The transformation will then add rules that charm additional terms by proving additional \( \text{magic}(\ldots) \) items (known as magic templates). Informally, a term needs to be charmed if it might contribute to the value of another charmed term. A formal description appears in Figure 7.

Finally, filter clauses are added to say that among the ground terms provable under the original program, only the charmed ones actually need to be proved. This means in practice (see section 6.5.2) that forward chaining will only prove an item if at least one grounding of that item is charmed.

The filter clauses in the speculation transformation were effectively introduced by explicitly charming all non-slashed items, running the magic templates transformation, and simplifying the result.

### 6.6.2 Examples of Magic Templates

**Deriving Earley’s algorithm.** What happens if we apply this transformation to the CKY algorithm of section 6.2.1, after explicitly charming the items shown above?

Remarkably, as previously noticed by Minnen (1996), the transformed program acts just like Earley’s (1970) algorithm. We can derive a weighted version of Earley’s algorithm by beginning with a weighted version of CKY (the inside algorithm of section 6.2.2). The transformation adds filter clauses to the \( \text{constit} \) rules, saying that the rule’s head is needed only if charmed:

\[
\begin{align*}
\text{constit}(X,I,K) & \rightarrow \text{rewrite}(X,W) \ast \text{word}(W,I,K) \\
& \quad \text{needed only if} \; \text{magic}(\text{constit}(X,I,K)). \\
\text{constit}(X,I,K) & \rightarrow \text{rewrite}(X,Y,Z) \ast \text{constit}(Y,I,J) \ast \text{constit}(Z,J,K) \\
& \quad \text{needed only if} \; \text{magic}(\text{constit}(X,I,K)).
\end{align*}
\]

\[44\] This terminology does not appear in previous literature on magic templates.

\[45\] At least, Earley’s algorithm restricted to grammars in Chomsky Normal Form, since those are the only grammars that CKY handles before we transform it. The full Earley’s algorithm in roughly our notation can be found in (Eisner et al., 2005); it allows arbitrary CFG rules to be expressed using lists, as in \( \text{rewrite}(\text{np},[\text{the},\text{adj},n]) \). Section 6.3 already sketched how to handle ternary rules efficiently.
Based on the structure of the \texttt{constit} rules, the transformation also adds the following magic rules to the ones provided earlier by the user. Recall that \texttt{?rewrite(X,Y,Z)} is considered to be true iff \texttt{rewrite(X,Y,Z)} is provable (footnote 12).

\begin{align*}
\text{magic(\texttt{constit(s,0,N)})} & = \text{magic(goal)}, \\
\text{magic(\texttt{constit(Y,I,J)})} & = \text{?rewrite(X,Y,Z) \& magic(\texttt{constit(X,I,K)})}, \\
\text{magic(\texttt{constit(Z,J,K)})} & = \text{?rewrite(X,Y,Z) \& ?constit(Y,I,J) \& magic(\texttt{constit(X,I,K)})}.
\end{align*}

What do these rules mean? The second magic rule above says to charm all the possible left children \texttt{constit(Y,I,J)} of a charmed constituent \texttt{constit(X,I,K)}. The third magic rule says to charm all possible right children of a charmed constituent whose left child has already been proved.

By inspecting these rules, one can see inductively that they prove only magic templates of the form \texttt{magic(\texttt{constit(X,I,K)})} where \texttt{X,I} are bound and \texttt{K} is free.

The charmed constituents are exactly those constituents that are possible given the context to their left. As in Earley’s algorithm, these are the only ones that we try to build. Where Earley’s algorithm would predict a \texttt{vp} constituent starting at position 5, we charm all potential constituents of that form by proving the non-ground item \texttt{magic(\texttt{constit(vp,5,K)})}.

Just as Earley’s predictions occur top-down, the magic rules reverse the proof order of the original rule—they charm items in the body of the original rule once the head is charmed. The magic rules also work left to right within an original rule, so we only need to charm \texttt{constit(vp,5,K)} once we have proved a receptive context such as \texttt{rewrite(s,np,vp) \& constit(np,4,5)}. This context is analogous to having the dotted rule \texttt{s \rightarrow np . vp} in column 5 of Earley’s parse table.

In effect, the transformed program uses forward chaining to simulate the backward-chaining proof order of a strategy like that of pure Prolog. The magic templates correspond to query subgoals that would appear during backward chaining. The filter clauses prevent the program from proving items that

\footnote{Note that it would not be appropriate to replace \ldots \& \texttt{?rewrite(X,Y,Z)} with \ldots \texttt{needed_only_if} \texttt{?rewrite(X,Y,Z)}, since that would make this condition optional, allowing the compiler to relax it and therefore charm more terms than intended. Concretely, in this example, forward chaining with our usual efficient treatment of \texttt{needed_only_if} (section 6.5.2) would prove overly general magic templates of the form \texttt{magic(\texttt{constit(X,I,K)})} where not only \texttt{K} but also \texttt{K} was free.}

\footnote{However, pure Prolog’s backtracking is deterministic, whereas forward chaining is free to propagate updates in any order. It is more precise to say that the transformed program simulates a \textit{breadth-first or parallel} version of Prolog which, when it has several ways to match a query subgoal, pursues them along parallel threads (whose actual operation on a serial computer may be interleaved in any way). Furthermore, since forward chaining uses a chart to avoid duplicate work, the transformed version acts like a version of Prolog with \textit{tabling} (see footnote 5).}
do not yet match any of these subgoals.

Shieber et al. (1995), specifying CKY and Earley’s algorithm, remark that “proofs of soundness and completeness [for the Earley’s case] are somewhat more complex . . . and are directly related to the corresponding proofs for Earley’s original algorithm.” In our perspective, the correctness of Earley’s emerges directly from the correctness of CKY and the correctness of the magic templates transformation (i.e., the fact that it preserves the semantics of charmed terms).

**On-demand computation of reachable states in a finite-state machine.** Another application is “on-the-fly” intersection of weighted finite-state automata, which recalls the left-to-right nature of Earley’s algorithm.48

In automaton intersection, an arc $Q_1 \xrightarrow{X} R_1$ (in some automaton $M_1$) may be paired with a similarly labeled arc $Q_2 \xrightarrow{X} R_2$ (in some automaton $M_2$, perhaps equal to $M_1$). This yields an arc in the intersected machine $M_1 \cap M_2$ whose weight is the product of the original arcs’ weights:

$$\text{arc}(Q_1:Q_2,R_1:R_2,X) = \text{arc}(Q_1,R_1,X) \times \text{arc}(Q_2,R_2,X).$$

However, including this rule in a forward-chained program will pairs all compatible arcs in all known machines (including arcs in the new machine $M_1 \cap M_2$, leading to infinite regress). A magic templates transformation can restrict this to arcs that actually need to be derived in the service of some larger goal—just as if the definition above were backward-chained.

Consider for example the following useful program (which uses Prolog’s notation for bracketed lists):

```
sum(Q,[],) += final(Q). % weight of stopping at final state Q
sum(Q,[X|Xs]) += arc(Q,R,X) \times sum(R,Xs).
```

Now the value of $\text{sum}(q,["a","b","c"])$ is the total weight of all paths from state $q$ that accept the string $abc$.

We might like to find (for example) $\text{sum}(q1;q2,["a","b","c"])$, constructing just the necessary paths from state $q1;q2$ in the intersection of $q1$’s automaton with $q2$’s automaton. To enable such queries, apply magic templates transformation to the $\text{sum}$ rules and the arc intersection rule, charming nothing in advance. We can then set $\text{magic(sum(q1;q2,["a","b","c"]))}$ to true at runtime. This results in charm spreading forward from $q1;q2$ along paths in the intersected machine, and spreading “top-down” from each arc along this path to the arcs that must be intersected to produce it (and which may themselves be the result of intersections). This permits the weights of all relevant arcs to be computed “bottom-up” and multiplied together along the desired paths.

48Composition of finite-state transducers is similar.
6.6.3 Second-order magic

Earley’s algorithm does top-down prediction quite aggressively, since it predicts all possible constituents that are consistent with the left context. Many of these predictions could be ruled out with one word of lookahead—an important technique when using Earley’s algorithm in practice. This is known as a “left-corner” filter: we should only bother to prove magic(constit(X,I,K)) if there is some chance of proving constit(X,I,K), in the sense that there is a word(W,I,J) that could serve as the first word (“left corner”) in a phrase constit(X,I,K).

Remarkably, we can get this behavior automatically by applying the magic templates transformation a second time. We now require the magic items themselves to be charmed before we will derive them. This activation flows bottom-up in the parse tree: we first charm magic(constit(X,I,K)) where X can rewrite directly as W, then move up to nonterminals that can rewrite starting with X, and so on.

Thus, the original CKY algorithm proved constituents bottom-up; the transformed Earley’s algorithm filtered these constituents by top-down predictions; and the doubly transformed algorithm will filter the predictions by bottom-up propagation of left corners.

Before illustrating this transformation, we point out some simplifications that are possible with second-order magic. This time we will use magic2 to indicate charmed items, to avoid conflict with the magic predicate that already appears in the input program. We also assume that the user is willing to explicitly charm everything but the magic terms—since any other terms that the user regards as uninteresting are presumably already being filtered by the last transformation. Suppose that the original program contained the rule a += b * c. The input program then also usually contains

\[
\begin{align*}
\text{magic}(b) & \Rightarrow \text{magic}(a), \\
\text{magic}(c) & \Rightarrow ?b \& \text{magic}(a).
\end{align*}
\]

However, either of these rules may be omitted if the user explicitly charmed its head during the first round of magic (i.e., by stating magic(b) = true or magic(c) = true). As we will see, omitting these rules when possible will reduce the work that second-order magic has to do.

If we apply a second round of magic literally, the above rules (when present) respectively yield the new rules

\[
\begin{align*}
\text{magic2}(\text{magic}(a)) & \Rightarrow \text{magic2}(\text{magic}(b)), \\
\text{magic2}(\text{magic}(b)) & \Rightarrow \text{magic2}(\text{magic}(c)), \\
\text{magic2}(\text{magic}(a)) & \Rightarrow ?b \& \text{magic2}(\text{magic}(c)).
\end{align*}
\]

These rules propagate charm on the magic items, from magic(b) or magic(c)

---

Earley (1970) himself described how to use k words of lookahead.
up to magic(a). However, it turns out that the second and often the third of these magic2 rules can be discarded, as they are redundant with more lenient rules that prove the same heads. The second is redundant because the user has already explicitly stated that magic2(b) |= true. The third is redundant if the first is present, since if the program has proved b then it must have previously proved magic(b) and before that magic2(magic(b)), so that the first rule (if present) would already have been able to prove magic2(magic(a)).

The input program also contains
\[ a + b \cdot c \text{ needed only if } magic(a). \]
We want to prove magic(a) if it will be useful in such a clause, so the second round of magic will also generate
\[ magic2(magic(a)) |= \exists b \land \exists c \land magic2(a). \]
This is the rule that initiates the desired bottom-up charming of magic items. It too can be simplified. We can drop the magic2(a) condition, since the user has already explicitly stated that magic2(a) |= true. We can drop the \( ?c \) condition if the third magic2 rule above is present, and drop the entire rule if the first magic2 rule above is present. (Thus, we will end up discarding the rule unless b was charmed by the user prior to the first round of magic—making it the appropriate “bottom” where bottom-up propagation begins.)

Applying second-order magic with these simplifications to our version of Earley’s algorithm, we obtain the following natural rules for introducing and propagating left corners. Note that these affect only the constit terms. Intuitively, the other terms of the original program do not need magic2 templates to entitle them to acquire first-order charm, as they were explicitly charmed by the user prior to first-order magic.

\[
\begin{align*}
t & \text{magic2(magic(constit(X,I,K)))) } \mid = \text{rewrite}(X,W) \land \text{word}(W,I,K), \\
t & \text{magic2(magic(constit(X,I,K)))) } \mid = \exists \text{rewrite}(X,Y,Z) \land \text{magic2(magic(constit(Y,I,J)).)}
\end{align*}
\]

The transformation then applies the left-corner filter to the magic templates defined by first-order magic:

\[
\begin{align*}
t & \text{magic(constit(s,0,N)) } \mid = \text{magic(goal) } \text{ needed only if } magic2(magic(constit(s,0,N))). \\
t & \text{magic(constit(Y,I,J)) } \mid = \exists \text{rewrite}(X,Y,Z) \land \text{magic(constit(X,I,K))} \land \text{magic2(magic(constit(Y,I,J)).)} \\
t & \text{magic(constit(Z,J,K)) } \mid = \exists \text{rewrite}(X,Y,Z) \land \exists \text{constit}(Y,I,J) \land \text{magic(constit(X,I,K)).} \text{ needed only if } magic2(magic(constit(Z,J,K)).)}
\end{align*}
\]

Note that the magic2(constit(X,I,K)) items proved above are specific to the span I–K in the current sentence: they have X,I,K all bound. However, one could remove this dependence on I,K by using the speculation transformation (section 6.5). Then the first time a particular word W is observed via some fact
Given a unary predicate magic that may already appear in \( P \). We say that a term \( t \) is already charmed if \( P \) contains a rule \( \text{magic}(s) |\!| \text{true} \) where \( s \) is at least as general as \( t \).

For each rule \( R_i \) in \( P \), of the form \( r_i \oplus E_i \), given an ordering \( e_1, \ldots, e_k \) of the items whose values are combined by \( E_i \) (including any filter clauses).\(^{51}\)

\[
\text{foreach rule } R_i
\]
\[
\quad \text{unless } r_i \text{ is already charmed}
\quad \text{append “needed only if } \text{magic}(r_i) \text{” to } R_i
\]
\[
\quad \text{for } j = 1, 2, \ldots, k_i
\quad \quad \text{unless } e_{ij} \text{ is already charmed}
\quad \quad \quad \text{add “} \text{magic}(e_{ij}) | \!| ? \text{e}_{ij} \& \cdots \& ? \text{e}_{i(j-1)} \& \text{magic}(r_i) \text{” to } P
\quad \quad \text{optionally relax this new rule by generalizing its head}^{52}\]

\(^{50}\)This test is used only to simplify the output.

\(^{51}\)In the examples in the text, this is taken to be the order of mention, which is a reasonable default.

\(^{52}\)That is, replace the head with a more general pattern. For example, one may replace some variables or other sub-terms in the head with variables that do not appear in the rule body. See section 6.6.4 for discussion.

**FIGURE 7** The magic templates transformation.

word(\(W,I,K\)), the program will derive magic2\((\text{constit}(X,I0,K0))/\text{word}(W,I0,K0)\) for each nonterminal \(X\) of which \(W\) is a possible left corner. These left corner table entries leave \(I0,K0\) free, so once they are computed, they can be combined not only with \(\text{word}(W,I,K)\) but also with later appearances of the same word, such as \(\text{word}(W,I2,K2)\).

### 6.6.4 Formalizing Magic Templates

Our version of magic templates is shown in Figure 7. Readers who are familiar with Ramakrishnan (1991) should note that our presentation focuses on the case that Ramakrishnan calls “full sips,” where each term used in a rule’s body constrains the bindings of variables in subsequent terms.

However, to allow other “sips” (sideways information-passing strategies), we can optionally rename variables in the heads of \(\text{magic}(\ldots) | \!| \cdots \text{ rules so that they become free. This results in proving fewer, more general } \text{magic}(\ldots) \text{ items.}^{55}\)

Ramakrishnan’s construction instead attempts to drop these variables—as well as other variables that provably remain free. However, his construction \(^{53}\)This may even avoid an asymptotic slowdown. Why? It is possible to prove more magic templates than items in the original program, because \(\text{magic}(a)\) (proved top-down) may acquire bindings for variables that are still free in \(a\) (proved bottom-up). It is wise to drop such variables from \(\text{magic}(a)\) or leave them free.
is less flexible because it only drops variables that appear as direct arguments to the top-level predicate. It also leads to a proliferation of new and non-interacting predicates (such as magic,constit/*), which correspond to different binding patterns in the top-level predicate.

Dropping variables rather than freeing them does have the advantage that it makes terms smaller, perhaps resulting in constant-time reductions of speed and space. However, we opt to defer this kind of “abbreviation” of terms to an optional subsequent transformation—the same transformation (not given in this paper) that we would use to abbreviate the slash(...) items introduced by speculation in section 6.5.1. Pushing abbreviation into a separate transformation keeps our Figure 7 simple. The abbreviation transformation could also attempt more ambitious simplifications than Ramakrishnan (1991), who does not simplify away nested free variables, duplicated bound variables, or constants, nor even detect all free variables that are arguments to the top-level predicate.

6.7 Conclusions

This paper has introduced the formalism of weighted logic programming, a powerful declarative language for describing a wide range of useful algorithms.

We outlined several fundamental techniques for rearranging a weighted logic program to make it more efficient. Each of the techniques is connected to ideas in both logic programming and in parsing, and has multiple uses in natural language processing. We used them to recovered several known parsing optimizations, such as

- unary rule cycle elimination
- Earley’s (1970) algorithm with an added left corner filter
- Eisner and Satta’s (1999) $O(n^3)$ bilexical parsing
- on-the-fly intersection of weighted automata

as well as various other small rearrangements of algorithms, such as a slight improvement to lexicalized CKY parsing.

We showed how weighted logic programming can be made more expressive and its transformations simplified by allowing non-ground items to be derived, and we introduced a new kind of side condition that does not bind variables—the needed_only_if construction—to streamline the use of non-ground items.

Our specific techniques included weighted generalizations of folding and unfolding; the speculation transformation (an original generalization of folding); and an improved formulation of the magic templates transformation. This work does not exhaust the set of useful transformations. For example,
Eisner et al. (2005) briefly discuss transformations that derive programs to compute gradients of, or bounds on, the values computed by the original dynamic program. We intend in the future to give formal treatments of these. We also plan to investigate other potentially useful transformations, in particular, transformations that exploit program invariants to perform tasks such as “abbreviating” complex items.

We hope that the paradigm presented here proves useful to those who wish to further study the problems to which weighted logic programming can be applied, as well as to those who wish to apply it to those problems themselves.

In the long run, we hope that by detailing a set of possible program transformation steps, we can work toward creating a system that would search automatically for practically effective transformations of a given weighted logic program, by incorporating observations about the program’s structure as well as data collected from experimental runs. Such an implemented system could be of great practical value.

References


On theoretical and practical complexity of TAG parsers

CARLOS GÓMEZ-RODRÍGUEZ, MIGUEL A. ALONSO, MANUEL VILARES

Abstract
We present a system allowing the automatic transformation of parsing schemata to efficient executable implementations of their corresponding algorithms. This system can be used to easily prototype, test and compare different parsing algorithms. In this work, it has been used to generate several different parsers for Context Free Grammars and Tree Adjoining Grammars. By comparing their performance on different sized, artificially generated grammars, we can measure their empirical computational complexity. This allows us to evaluate the overhead caused by using Tree Adjoining Grammars to parse context-free languages, and the influence of string and grammar size on Tree Adjoining Grammars parsing.

Keywords Parsing schemata, computational complexity, tree adjoining grammars, context free grammars

7.1 Introduction
The process of parsing, by which we obtain the structure of a sentence as a result of the application of grammatical rules, is a highly relevant step in the automatic analysis of natural languages. In the last decades, various parsing algorithms have been developed to accomplish this task. Although all of these algorithms essentially share the common goal of generating a tree structure describing the input sentence by means of a grammar, the approaches used to attain this result vary greatly between algorithms, so that different parsing algorithms are best suited to different situations.

Parsing schemata, introduced in (Sikkel, 1997), provide a formal, simple and uniform way to describe, analyze and compare different parsing algorithms. The notion of a parsing schema comes from considering parsing as a
deduction process which generates intermediate results called items. An initial set of items is directly obtained from the input sentence, and the parsing process consists of the application of inference rules (called deductive steps) which produce new items from existing ones. Each item contains a piece of information about the sentence’s structure, and a successful parsing process will produce at least one final item containing a full parse tree for the sentence or guaranteeing its existence.

Almost all known parsing algorithms may be described by a parsing schema (non-constructive parsers, such as those based on neural networks, are exceptions). This is done by identifying the kinds of items that are used by a given algorithm, defining a set of inference rules describing the legal ways of obtaining new items, and specifying the set of final items.

As an example, we introduce a CYK-based algorithm (Vijay-Shanker and Joshi 1985) for Tree Adjoining Grammars (TAG) (Joshi and Schabes 1997). Given a tree adjoining grammar $G = (V_T, V_N, S, I, A)$\footnote{Where $V_T$ denotes the set of terminal symbols, $V_N$ the set of nonterminal symbols, $S$ the axiom, $I$ the set of initial trees and $A$ the set of auxiliary trees.} and a sentence of length $n$ which we denote by $a_1 a_2 ... a_n$, we denote by $P(G)$ the set of productions $\{N_\gamma \rightarrow N_{\gamma 1} N_{\gamma 2} ... N_{\gamma r}\}$ such that $N_\gamma$ is an inner node of a tree $\gamma \in (I \cup A)$, and $N_{\gamma 1} N_{\gamma 2} ... N_{\gamma r}$ is the ordered sequence of direct children of $N_\gamma$.

The parsing schema for the TAG CYK-based algorithm is a function that maps such a grammar $G$ to a deduction system whose domain is the set of items $\{[N_\gamma, i, j, p, q, adj]\}$ verifying that $N_\gamma$ is a tree node in an elementary tree $\gamma \in (I \cup A)$, $i$ and $j$ ($0 \leq i \leq j$) are string positions, $p$ and $q$ may be undefined or instantiated to positions $i \leq p \leq q \leq j$ (the latter only when $\gamma \in A$), and $adj \in \{true, false\}$ indicates whether an adjunction has been performed on node $N_\gamma$.

The positions $i$ and $j$ indicate that a substring $a_{i+1} ... a_{j}$ of the string is being recognized, and positions $p$ and $q$ denote the substring dominated by $\gamma$’s foot node. The final item set would be $\{[R_\alpha, 0, n, -, -, adj] | \alpha \in I\}$ for the presence of such an item would indicate that there exists a valid parse tree with yield $a_1 a_2 ... a_n$ and rooted at $R_\alpha$, the root of an initial tree; and therefore there exists a complete parse tree for the sentence.

A deductive step $\frac{[R_\alpha]}{\frac{\xi}{\Phi}}$ allows us to infer the item specified by its con-
sequent $\xi$ from those in its antecedents $\eta_1 \ldots \eta_m$. Side conditions ($\Phi$) specify the valid values for the variables appearing in the antecedents and consequent, and may refer to grammar rules or specify other constraints that must be verified in order to infer the consequent. An example of one of the schema’s deductive steps would be the following, where the operation $p \cup p'$ returns $p$ if $p$ is defined, and $p'$ otherwise:

$$
\begin{align*}
\text{CYK Binary:} \\
[O_{\gamma 1}, i, j, p, q, \text{adj} 1] \\
[O_{\gamma 2}, j', j, p', q', \text{adj} 2] \\
[M', i, j, p \cup p', q \cup q', false] \\
M' \rightarrow O_{\gamma 1} O_{\gamma 2} \in \text{P}(G)
\end{align*}
$$

This deductive step represents the bottom-up parsing operation which joins two subtrees into one, and is analogous to one of the deductive steps of the CYK parser for Context-Free Grammars (Kasami 1965, Younger 1967). The full TAG CYK parsing schema has six deductive steps (or seven, if we work with TAGs supporting the substitution operation) and can be found at (Alonso et al., 1999). However, this sample deductive step is an example of how parsing schemata convey the fundamental semantics of parsing algorithms in simple, high-level descriptions. A parsing schema defines a set of possible intermediate results and allowed operations on them, but doesn’t specify data structures for storing the results or an order for the operations to be executed.

### 7.2 Compilation of parsing schemata

Their simplicity and abstraction of low-level details makes parsing schemata very useful, allowing us to define parsers in a simple and straightforward way. Comparing parsers, or considering aspects such as their correction and completeness or their computational complexity, also becomes easier if we think in terms of schemata.

However, the problem with parsing schemata is that, although they are very useful when designing and comparing parsers with pencil and paper, they cannot be executed directly in a computer. In order to execute the parsers and analyze their results and performance they must be implemented in a programming language, making it necessary to abandon the high abstraction level and focus on the implementation details in order to obtain a functional and efficient implementation.

In order to bridge this gap between theory and practice, we have designed and implemented a compiler able to automatically transform parsing schemata into efficient Java implementations of their corresponding algorithms. The input to this system is a simple and declarative representation of a parsing schema, which is practically equal to the formal notation that we used previously. For example, this is the CYK deductive step we have seen as an example in a format readable by our compiler:
The parsing schema compilation technique behind our system is based on the following fundamental ideas:

- Each deductive step is compiled to a Java class containing code to match and search for antecedent items and generate the corresponding conclusions from the consequent.

- The generated implementation will create an instance of this class for each possible set of values satisfying the side conditions that refer to production rules. For example, a distinct instance of the CYK Binary step will be created for each grammar rule of the form \( M^\gamma \rightarrow O^\gamma_1 O^\gamma_2 \in P(G) \), as specified in the step’s side condition.

- The step instances are coordinated by a deductive parsing engine, as the one described in (Shieber et al., 1995). This algorithm ensures a sound and complete deduction process, guaranteeing that all items that can be generated from the initial items will be obtained. It is a generic, schema-independent algorithm, so its implementation is the same for any parsing schema. The engine works with the set of all items that have been generated and an agenda, implemented as a queue, holding the items we have not yet tried to trigger new deductions with.

- In order to attain efficiency, an automatic analysis of the schema is performed in order to create indexes allowing fast access to items. Two kinds of index structures are generated: existence indexes are used by the parsing engine to check whether a given item exists in the item set, while search indexes are used to search for all items conforming to a given specification. As each different parsing schema needs to perform different searches for antecedent items, the index structures that we generate are schema-specific. Each deductive step is analyzed in order to keep track of which variables will be instantiated to a concrete value when a search must be performed. This information is known at schema compilation time and allows us to create indexes by the elements corresponding to instantiated variables. In this way, we guarantee constant-time access to items so that the computational complexity of our generated implementations is never above the theoretical complexity of the parsing algorithms.

- Deductive step indexes are also generated to provide efficient access to the set of deductive step instances which can be applicable to a given item. Step instances that are known not to match the item are filtered out by
these indexes, so less time is spent on unsuccessful item matching.

- Since parsing schemata have an open notation, for any mathematical object can potentially appear inside items, the system includes an extensibility mechanism which can be used to define new kinds of objects to use in schemata. The code generator can deal with these user-defined objects as long as some simple and well-defined guidelines are followed in their specification.

A more detailed description of this system, including a more thorough explanation of automatic index generation, can be found at (Gómez-Rodríguez et al., 2006b, 2007).

7.3 Parsing natural language CFGs

Although our main focus in this paper is on performance of TAG parsing algorithms, we will briefly outline the results of some experiments on Context-Free Grammars (CFG), described in further detail in (Gómez-Rodríguez et al., 2006b), in order to be able to contrast TAG and CFG parsing.

Our compilation technique was used to generate parsers for the CYK (Kasami 1965, Younger 1967), Earley (Earley 1970) and Left-Corner (Rosenkrantz and Lewis II 1970) algorithms for context-free grammars, and these parsers were tested on automatically-generated sentences from three different natural language grammars: Susanne (Sampson 1994), Alvey (Carroll 1993) and Deltra (Schoorl and Belder 1990). The run-times for all the algorithms and grammars showed an empirical computational complexity far below the theoretical worst-case bound of \( O(n^3) \), where \( n \) denotes the length of the input string. In the case of the Susanne grammar, the measurements were close to being linear with string size. In the other grammars, the run-times grew faster, approximately \( O(n^2) \), still far below the cubic worst-case bound.

Another interesting result was that the CYK algorithm performed better than the Earley-type algorithms in all cases, despite being generally considered slower. The reason is that these considerations are based on time complexity relative to string length, and do not take into account time complexity relative to grammar size, which is \( O(|P|) \) for CYK and \( O(|P|^2) \) for the Earley-type algorithms, where \( |P| \) is the number of production rules in the grammar. This factor is not very important when working with small grammars, such as the ones used for programming languages, but it becomes fundamental when we work with natural language grammars, where we use thousands of rules (more than 17,000 in the case of Susanne) to parse relatively small sentences. When comparing the results from the three context-free grammars, we observed that the performance gap between CYK and Earley was bigger when working with larger grammars.\(^3\)

\(^3\)It is possible to reduce the computational complexity of Earley’s parser to linear with respect
7.4 Parsing artificial TAGs

In this section, we make a comparison of four different TAG parsing algorithms: the CYK-based algorithm used as an example in section 7.1, an Earley-based algorithm without the valid prefix property (described in Alonso et al. 1999 and Alonso et al. 2004, inspired in the one in Schabes 1994), an Earley-based algorithm with the valid prefix property (Alonso et al. 1999) and Nederhof’s algorithm (Nederhof 1999, Alonso et al. 2004). These parsers are compared on artificially generated grammars, by using our schema compiler to generate implementations and measuring their execution times with several grammars and sentences.

Note that the advantage of using artificially generated grammars is that we can easily see the influence of grammar size on performance. If we test the algorithms on grammars from real-life natural language corpora, as we did with the CFG parsers, we don’t get a very precise idea of how the size of the grammar affects performance. Since our experience with CFGs showed this to be an important factor, and existing TAG parser performance comparisons (e.g. Díaz and Alonso 2000) work with a fixed (and small) grammar, we decided to use artificial grammars in order to be able to adjust both string size and grammar size in our experiments and see the influence of both factors.

For this purpose, given an integer \( k > 0 \), we define the tree-adjoining grammar \( G_k \) to be the grammar \( G_k = (V_T, V_N, S, I, A) \) where \( V_T = \{ a_j | 0 \leq j \leq k \} \), \( V_N = \{ S, B \} \), and

\[
I = \{ S(B(a_0)) \}^k, \\
A = \{ B(B^* a_j) \} | 1 \leq j \leq k \}.
\]

Therefore, for a given \( k \), \( G_k \) is a grammar with one initial tree and \( k \) auxiliary trees, which parses a language over an alphabet with \( k + 1 \) terminal symbols. The actual language defined by \( G_k \) is the regular language \( L_k = a_0(a_1|a_2|...|a_k)^* \). We shall note that although the languages \( L_k \) are trivial, the grammars \( G_k \) are built in such a way that any of the auxiliary trees may adjoin into any other. Therefore these grammars are suitable if we want to make an empirical analysis of worst-case complexity.

to the grammar size by defining a new set of intermediate items and transforming accordingly prediction and completion deduction steps. Even in this case, CYK performs better that Earley’s algorithm due to the lower number of items generated: \( O(|V_N| + |V_T| n^2) \) for CYK vs. \( O(|G| n^2) \) for Earley’s algorithm, where \(|G|\) denotes the size of the grammar measured as \(|P|\) plus the summation of the lengths of all productions.

Where trees are written in bracketed notation, and * is used to denote the foot node.

Also, it is easy to prove that the grammar \( G_k \) is one of the minimal tree adjoining grammars (in terms of number of trees) whose associated language is \( L_k \). Note that we need at least a tree containing \( a_0 \) as its only terminal in order to parse the sentence \( a_0 \), and for each \( 1 \leq i \leq k \), we need at least a tree containing \( a_i \) and no other \( a_j \) \( (j > 0) \) in order to parse the sentence \( a_0a_j \). Therefore, any TAG for the language \( L_k \) must have at least \( k + 1 \) elementary trees.
Table 1 shows the execution time in milliseconds⁶ of four TAG parsers with the grammars $G_k$, for different values of string length ($n$) and grammar size ($k$).

From this results, we can observe that both factors (string length and grammar size) have an influence on runtime, and they interact between themselves: the growth rates with respect to one factor are influenced by the other factor, so it is hard to give precise estimates of empirical computational complexity. However, we can get rough estimates by focusing on cases where one of the factors takes high values and the other one takes low values (since in these cases the constant factors affecting complexity will be smaller) and test them by checking whether the sequence $T(n,k)/f(n)$ seems to converge to a positive constant for each fixed $k$ (if $f(n)$ is an estimation of complexity with respect to string length) or whether $T(n,k)/f(k)$ seems to converge to a positive constant for each fixed $n$ (if $f(k)$ is an estimation of complexity with respect to grammar size).

By applying these principles, we find that the empirical time complexity with respect to string length is in the range between $O(n^{2.8})$ and $O(n^3)$ for the CYK-based and Nederhof algorithms, and between $O(n^{2.6})$ and $O(n^3)$ for the Earley-based algorithms with and without the valid prefix property (VPP). Therefore, the practical time complexity we obtain is far below the theoretical worst-case bounds for these algorithms, which are $O(n^6)$ (except for the Earley-based algorithm with the VPP, which is $O(n^7)$).

Although for space reasons we don’t include tables with the number of items generated in each case, our results show that the empirical space complexity with respect to string length is approximately $O(n^2)$ for all the algorithms, also far below the worst-case bounds ($O(n^4)$ and $O(n^5)$).

With respect to the size of the grammar, we obtain a time complexity of approximately $O(|I∪A|^2)$ for all the algorithms. This matches the theoretical worst-case bound, which is $O(|I∪A|^2)$ due to the adjunction steps, which work with pairs of trees. In the case of our artificially generated grammar, any auxiliary tree can adjoin into any other, so it’s logical that our times grow quadratically. Note, however, that real-life grammars such as the XTAG English grammar (XTAG Research Group 2001) have relatively few different nonterminals in relation to their amount of trees, so many pairs of trees are susceptible of adjunction and we can’t expect their behavior to be much better than this.

Space complexity with respect to grammar size is approximately $O(|I∪A|)$ for all the algorithms. This is an expected result, since each generated item is associated to a given tree node.

---

⁶The machine used for all the tests was an Intel Pentium 4 3.40 GHz, with 1 GB RAM and Sun Java Hotspot virtual machine (version 1.4.2_01-b06) running on Windows XP.
<table>
<thead>
<tr>
<th>String Size (n)</th>
<th>1</th>
<th>8</th>
<th>64</th>
<th>512</th>
<th>4096</th>
</tr>
</thead>
<tbody>
<tr>
<td>2</td>
<td>-0</td>
<td>16</td>
<td>15</td>
<td>1,156</td>
<td>109,843</td>
</tr>
<tr>
<td>4</td>
<td>-0</td>
<td>31</td>
<td>63</td>
<td>2,578</td>
<td>256,094</td>
</tr>
<tr>
<td>8</td>
<td>16</td>
<td>31</td>
<td>172</td>
<td>6,891</td>
<td>589,578</td>
</tr>
<tr>
<td>16</td>
<td>31</td>
<td>172</td>
<td>625</td>
<td>18,735</td>
<td>1,508,609</td>
</tr>
<tr>
<td>32</td>
<td>110</td>
<td>609</td>
<td>3,219</td>
<td>69,406</td>
<td></td>
</tr>
<tr>
<td>64</td>
<td>485</td>
<td>2,953</td>
<td>22,453</td>
<td>289,984</td>
<td></td>
</tr>
<tr>
<td>128</td>
<td>2,031</td>
<td>13,875</td>
<td>234,594</td>
<td></td>
<td></td>
</tr>
<tr>
<td>256</td>
<td>10,000</td>
<td>101,219</td>
<td></td>
<td></td>
<td></td>
</tr>
<tr>
<td>512</td>
<td>61,266</td>
<td></td>
<td></td>
<td></td>
<td></td>
</tr>
</tbody>
</table>

**Run-times in ms: CYK-based**

<table>
<thead>
<tr>
<th>String Size (n)</th>
<th>1</th>
<th>8</th>
<th>64</th>
<th>512</th>
<th>4096</th>
</tr>
</thead>
<tbody>
<tr>
<td>2</td>
<td>-0</td>
<td>-0</td>
<td>16</td>
<td>1,344</td>
<td>125,750</td>
</tr>
<tr>
<td>4</td>
<td>-0</td>
<td>-0</td>
<td>63</td>
<td>4,109</td>
<td>290,187</td>
</tr>
<tr>
<td>8</td>
<td>16</td>
<td>31</td>
<td>234</td>
<td>15,891</td>
<td>777,968</td>
</tr>
<tr>
<td>16</td>
<td>15</td>
<td>62</td>
<td>782</td>
<td>44,188</td>
<td>2,247,156</td>
</tr>
<tr>
<td>32</td>
<td>94</td>
<td>312</td>
<td>3,781</td>
<td>170,609</td>
<td></td>
</tr>
<tr>
<td>64</td>
<td>266</td>
<td>2,063</td>
<td>25,094</td>
<td>550,016</td>
<td></td>
</tr>
<tr>
<td>128</td>
<td>1,187</td>
<td>14,516</td>
<td>269,047</td>
<td></td>
<td></td>
</tr>
<tr>
<td>256</td>
<td>6,781</td>
<td>108,297</td>
<td></td>
<td></td>
<td></td>
</tr>
<tr>
<td>512</td>
<td>52,000</td>
<td></td>
<td></td>
<td></td>
<td></td>
</tr>
</tbody>
</table>

**Run-times in ms: Nederhof’s Algorithm**

<table>
<thead>
<tr>
<th>String Size (n)</th>
<th>1</th>
<th>8</th>
<th>64</th>
<th>512</th>
<th>4096</th>
</tr>
</thead>
<tbody>
<tr>
<td>2</td>
<td>-0</td>
<td>-0</td>
<td>47</td>
<td>1,875</td>
<td>151,532</td>
</tr>
<tr>
<td>4</td>
<td>-0</td>
<td>15</td>
<td>187</td>
<td>4,563</td>
<td>390,468</td>
</tr>
<tr>
<td>8</td>
<td>15</td>
<td>31</td>
<td>469</td>
<td>12,531</td>
<td>998,594</td>
</tr>
<tr>
<td>16</td>
<td>46</td>
<td>188</td>
<td>1,500</td>
<td>40,093</td>
<td>2,579,578</td>
</tr>
<tr>
<td>32</td>
<td>219</td>
<td>953</td>
<td>6,235</td>
<td>157,063</td>
<td></td>
</tr>
<tr>
<td>64</td>
<td>1,078</td>
<td>4,735</td>
<td>35,860</td>
<td>620,047</td>
<td></td>
</tr>
<tr>
<td>128</td>
<td>5,703</td>
<td>25,703</td>
<td>302,766</td>
<td></td>
<td></td>
</tr>
<tr>
<td>256</td>
<td>37,125</td>
<td>159,609</td>
<td></td>
<td></td>
<td></td>
</tr>
<tr>
<td>512</td>
<td>291,141</td>
<td></td>
<td></td>
<td></td>
<td></td>
</tr>
</tbody>
</table>

**Run-times in ms: Earley-based with the VPP**

<table>
<thead>
<tr>
<th>String Size (n)</th>
<th>1</th>
<th>8</th>
<th>64</th>
<th>512</th>
<th>4096</th>
</tr>
</thead>
<tbody>
<tr>
<td>2</td>
<td>-0</td>
<td>-0</td>
<td>31</td>
<td>1,937</td>
<td>194,047</td>
</tr>
<tr>
<td>4</td>
<td>-0</td>
<td>16</td>
<td>78</td>
<td>4,078</td>
<td>453,203</td>
</tr>
<tr>
<td>8</td>
<td>15</td>
<td>31</td>
<td>234</td>
<td>10,922</td>
<td>781,141</td>
</tr>
<tr>
<td>16</td>
<td>31</td>
<td>188</td>
<td>875</td>
<td>27,125</td>
<td>1,787,140</td>
</tr>
<tr>
<td>32</td>
<td>125</td>
<td>750</td>
<td>4,141</td>
<td>98,829</td>
<td></td>
</tr>
<tr>
<td>64</td>
<td>578</td>
<td>3,547</td>
<td>28,640</td>
<td>350,218</td>
<td></td>
</tr>
<tr>
<td>128</td>
<td>2,453</td>
<td>20,766</td>
<td>264,500</td>
<td></td>
<td></td>
</tr>
<tr>
<td>256</td>
<td>12,187</td>
<td>122,797</td>
<td></td>
<td></td>
<td></td>
</tr>
<tr>
<td>512</td>
<td>74,046</td>
<td></td>
<td></td>
<td></td>
<td></td>
</tr>
</tbody>
</table>

**TABLE 1** Execution times of four different TAG parsers for artificially-generated grammars $G_k$. Best results are shown in boldface.
Practical applications of TAG in natural language processing usually fall in the range of values for $n$ and $k$ covered in our experiments (grammars with hundreds or a few thousands of trees are used to parse sentences of several dozens of words). Within these ranges, both string length and grammar size take significant values and have an important influence on execution times, as we can see from the results in the tables. This leads us to note that traditional complexity analysis based on a single factor (string length or grammar size) can be misleading for practical applications, since it can lead us to an incomplete idea of real complexity. For example, if we are working with a grammar with thousands of trees, the size of the grammar is the most influential factor, and the use of filtering techniques (Schabes and Joshi 1991) to reduce the amount of trees used in parsing is essential in order to achieve good performance. The influence of string length in these cases, on the other hand, is mitigated by the huge constant factors related to grammar size. For instance, in the times shown in the tables for the grammar $G_{1096}$, we can see that parsing times are multiplied by a factor less than 3 when the length of the input string is duplicated, although the rest of the results have lead us to conclude that the practical asymptotic complexity with respect to string length is at least $O(n^{2.6})$. These interactions between both factors must be taken into account when analyzing performance in terms of computational complexity.

Earley-based algorithms achieve better execution times than the CYK-based algorithm for large grammars, although they are worse for small grammars. This contrasts with the results for context-free grammars, where CYK works better for large grammars: when working with CFGs, CYK has a better computational complexity than Earley with respect to grammar size (see section 7.3), but the TAG variant of the CYK algorithm is quadratic with respect to grammar size and does not have this advantage.

CYK generates fewer items than the Earley-based algorithms when working with large grammars and short strings, and the opposite happens when working with small grammars and long strings.

The Earley-based algorithm with the VPP generates the same number of items than the one without this property, and has worse execution times. The reason is that no partial parses violating this property are generated by any of both algorithms in the particular case of this grammar, so guaranteeing the valid prefix property does not prevent any items from being generated. Therefore, the fact that the variant without the VPP works better in this particular case cannot be extrapolated to other grammars. However, the differences in times between these two algorithms illustrates the overhead caused by the extra checks needed to guarantee the valid prefix property in a particularly bad case.

Nederhof’s algorithm has slower execution times than the other Earley variants. Despite the fact that Nederhof’s algorithm is an improvement over
the other Earley-based algorithm with the VPP in terms of computational complexity, the extra deductive steps it contains makes it slower in practice.

7.5 Parsing the XTAG English grammar

In order to complement our performance comparison of the four algorithms on artificial grammars, we have also studied the behavior of the parsers when working with a real-life, large-scale TAG: the XTAG English grammar (XTAG Research Group 2001).

The obtained execution times are in the ranges that we could expect given the artificial grammar results, i.e. they approximately match the times in the tables for the corresponding grammar sizes and input string lengths. The most noticeable difference is that the Earley-like algorithm verifying the valid prefix property generates fewer items that the variant without the VPP in the XTAG grammar, and this causes its run-times to be faster. But this difference is not surprising, as explained in the previous section.

Note that, as the XTAG English grammar has over a thousand elementary trees, execution times are very large (over 100 seconds) when working with the full grammar, even with short sentences. However, when a tree selection filter is applied in order to work with only a subset of the grammar in function of the input string, the grammar size is reduced to one or two hundred trees and our parsers process short sentences in less than 5 seconds. Sarkar’s XTAG distribution parser written in C\(^7\) applies further filtering techniques and has specific optimizations for this grammar, obtaining better times for the XTAG than our generic parsers.

Table 2 contains a summary of the execution times obtained by our parsers for some sample sentences from the XTAG distribution. Note that the generated implementations used for these executions apply the mentioned tree filtering technique, so that the effective grammar size is different for each sentence, hence the high variability in execution times. More detailed information on these experiments with the XTAG English grammar can be found at (Gómez-Rodríguez et al., 2006a).

7.6 Overhead of TAG parsing over CFG parsing

The languages \(L_k\) that we parsed in section 7.4 were regular languages, so in practice we don’t need tree adjoining grammars to parse them, although it was convenient to use them in our comparison. This can lead us to wonder how large is the overhead caused by using the TAG formalism to parse context-free languages.

Given the regular language \(L_k = a_0(a_1|a_2|...|a_k)^*\), a context-free grammar that parses it is \(G_k' = (N, \Sigma, P, S)\) with \(N = \{S\}\) and

\(^7\)Downloadable at: ftp://ftp.cis.upenn.edu/pub/xtag/lem/
On theoretical and practical complexity of TAG parsers / 97

<table>
<thead>
<tr>
<th>Sentence</th>
<th>CYK</th>
<th>VPP</th>
<th>Ear.</th>
<th>VPP</th>
<th>Neder.</th>
</tr>
</thead>
<tbody>
<tr>
<td>He was a cow</td>
<td>2985</td>
<td>750</td>
<td>750</td>
<td></td>
<td>2719</td>
</tr>
<tr>
<td>He loved himself</td>
<td>3109</td>
<td>1562</td>
<td>2129</td>
<td>6421</td>
<td></td>
</tr>
<tr>
<td>Go to your room</td>
<td>4078</td>
<td>1547</td>
<td>1406</td>
<td>6828</td>
<td></td>
</tr>
<tr>
<td>He is a real man</td>
<td>4266</td>
<td>1563</td>
<td>1407</td>
<td>4703</td>
<td></td>
</tr>
<tr>
<td>He was a real man</td>
<td>4234</td>
<td>1921</td>
<td>1421</td>
<td>4766</td>
<td></td>
</tr>
<tr>
<td>Who was at the door</td>
<td>4485</td>
<td>1813</td>
<td>1562</td>
<td>7782</td>
<td></td>
</tr>
<tr>
<td>He loved all cows</td>
<td>5469</td>
<td>2359</td>
<td>2344</td>
<td>11469</td>
<td></td>
</tr>
<tr>
<td>He called up her</td>
<td>7828</td>
<td>4906</td>
<td>3563</td>
<td>15532</td>
<td></td>
</tr>
<tr>
<td>He wanted to go to the city</td>
<td>10047</td>
<td>4422</td>
<td>4016</td>
<td>18969</td>
<td></td>
</tr>
<tr>
<td>That woman in the city contributed to</td>
<td>13641</td>
<td>6515</td>
<td>7172</td>
<td>31828</td>
<td></td>
</tr>
<tr>
<td>That people are not really amateurs at</td>
<td>16500</td>
<td>7781</td>
<td>15235</td>
<td>56265</td>
<td></td>
</tr>
<tr>
<td>intellectual dueling</td>
<td></td>
<td></td>
<td></td>
<td></td>
<td></td>
</tr>
<tr>
<td>The index is intended to measure future</td>
<td>16875</td>
<td>17109</td>
<td>9985</td>
<td>39132</td>
<td></td>
</tr>
<tr>
<td>economic performance</td>
<td></td>
<td></td>
<td></td>
<td></td>
<td></td>
</tr>
<tr>
<td>They expect him to cut costs through-</td>
<td>25859</td>
<td>12000</td>
<td>20828</td>
<td>63641</td>
<td></td>
</tr>
<tr>
<td>out the organization</td>
<td></td>
<td></td>
<td></td>
<td></td>
<td></td>
</tr>
<tr>
<td>He will continue to place a huge burden</td>
<td>54578</td>
<td>35829</td>
<td>57422</td>
<td>178875</td>
<td></td>
</tr>
<tr>
<td>on the city workers</td>
<td></td>
<td></td>
<td></td>
<td></td>
<td></td>
</tr>
<tr>
<td>He could have been simply being a jerk</td>
<td>62157</td>
<td>113532</td>
<td>109062</td>
<td>133515</td>
<td></td>
</tr>
<tr>
<td>A few fast food outlets are giving it a</td>
<td>269187</td>
<td>3122860</td>
<td>3315359</td>
<td></td>
<td></td>
</tr>
<tr>
<td>try</td>
<td></td>
<td></td>
<td></td>
<td></td>
<td></td>
</tr>
</tbody>
</table>

Table 2: Run-times obtained by applying different XTAG parsers to several sentences. Best results for each sentence are shown in boldface.

\[ P = \{ S \rightarrow a_0 \} \cup \{ S \rightarrow S a_i | 1 \leq i \leq k \} \]

This grammar minimizes the number of rules needed to parse \( L_k \) (\( k + 1 \) rules), but has left recursion. If we want to eliminate left recursion, we can use the grammar \( G_k'' = (N, \Sigma, P, S) \) with \( N = \{ S, A \} \) and

\[ P = \{ S \rightarrow a_0 A \} \cup \{ A \rightarrow a_i A | 1 \leq i \leq k \} \cup \{ A \rightarrow \varepsilon \} \]

which has \( k + 2 \) production rules.

The number of items generated by the Earley algorithm for context-free grammars when parsing a sentence of length \( n \) from the language \( L_k \) by using the grammar \( G_k'' \) is \( (k + 2)n \). In the case of the grammar \( G_k''' \), the same algorithm generates \( (k + 4)n + \frac{(k+1)(k+2)}{2} + 1 \) items. In both cases the amount of items generated is linear with respect to grammar size, as in TAG parsers. With respect to string size, the amount of items is \( O(n) \) for \( G_k'' \) and \( O(n^2) \) for \( G_k''' \), and it was approximately \( O(n^2) \) for the TAG \( G_k \). Note, however, that the constant factors
behind complexity are much greater when working with \( G_k \) than with \( G''_k \), and this reflects on the actual number of items generated (for example, the Earley algorithm generates 16,833 items when working with \( G''_{64} \) and a string of length \( n = 128 \), while the TAG variant of Earley without the valid prefix property generated 1,152,834 items).

The execution times for both algorithms appear in table 3. From the obtained times, we can deduce that the empirical time complexity is linear with respect to string length and quadratic with respect to grammar size in the case of \( G'_{16} \), and quadratic with respect to string length and linear with respect to grammar size in the case of \( G''_{16} \). So this example shows that, when parsing a context-free language using a tree-adjoining grammar, we get an overhead both in constant factors (more complex items, more deductive steps, etc.) and in asymptotic behavior, so actual execution times can be several orders of magnitude larger. Note that the way grammars are designed also has an influence, but our tree adjoining grammars \( G_k \) are the simplest TAGs able to parse the languages \( L_k \) by using adjunction (an alternative would be to write a grammar using the substitution operation to combine trees).

<table>
<thead>
<tr>
<th>n</th>
<th>Grammar Size (k), grammar ( G'_{k} )</th>
<th>1</th>
<th>8</th>
<th>64</th>
<th>512</th>
<th>4096</th>
</tr>
</thead>
<tbody>
<tr>
<td>2</td>
<td>-0</td>
<td>-0</td>
<td>-0</td>
<td>-0</td>
<td>31</td>
<td>2,062</td>
</tr>
<tr>
<td>4</td>
<td>-0</td>
<td>-0</td>
<td>-0</td>
<td>-0</td>
<td>62</td>
<td>4,110</td>
</tr>
<tr>
<td>8</td>
<td>-0</td>
<td>-0</td>
<td>-0</td>
<td>-0</td>
<td>125</td>
<td>8,265</td>
</tr>
<tr>
<td>16</td>
<td>-0</td>
<td>-0</td>
<td>-0</td>
<td>-0</td>
<td>217</td>
<td>15,390</td>
</tr>
<tr>
<td>32</td>
<td>-0</td>
<td>-0</td>
<td>-0</td>
<td>15</td>
<td>563</td>
<td>29,344</td>
</tr>
<tr>
<td>64</td>
<td>-0</td>
<td>-0</td>
<td>31</td>
<td>1,062</td>
<td>61,875</td>
<td></td>
</tr>
<tr>
<td>128</td>
<td>-0</td>
<td>-0</td>
<td>109</td>
<td>2,083</td>
<td>122,875</td>
<td></td>
</tr>
<tr>
<td>256</td>
<td>-0</td>
<td>15</td>
<td>188</td>
<td>4,266</td>
<td>236,688</td>
<td></td>
</tr>
<tr>
<td>512</td>
<td>15</td>
<td>328</td>
<td>8,406</td>
<td>484,859</td>
<td></td>
<td></td>
</tr>
</tbody>
</table>

<table>
<thead>
<tr>
<th>n</th>
<th>Grammar Size (k), grammar ( G''_{k} )</th>
<th>1</th>
<th>8</th>
<th>64</th>
<th>512</th>
<th>4096</th>
</tr>
</thead>
<tbody>
<tr>
<td>2</td>
<td>-0</td>
<td>-0</td>
<td>-0</td>
<td>15</td>
<td>94</td>
<td></td>
</tr>
<tr>
<td>4</td>
<td>-0</td>
<td>-0</td>
<td>-0</td>
<td>16</td>
<td>203</td>
<td></td>
</tr>
<tr>
<td>8</td>
<td>-0</td>
<td>-0</td>
<td>-0</td>
<td>46</td>
<td>688</td>
<td></td>
</tr>
<tr>
<td>16</td>
<td>-0</td>
<td>-0</td>
<td>15</td>
<td>203</td>
<td>1,735</td>
<td></td>
</tr>
<tr>
<td>32</td>
<td>-0</td>
<td>-0</td>
<td>93</td>
<td>516</td>
<td>4,812</td>
<td></td>
</tr>
<tr>
<td>64</td>
<td>31</td>
<td>31</td>
<td>516</td>
<td>4,812</td>
<td></td>
<td></td>
</tr>
<tr>
<td>128</td>
<td>156</td>
<td>156</td>
<td>1,500</td>
<td>13,406</td>
<td></td>
<td></td>
</tr>
<tr>
<td>256</td>
<td>484</td>
<td>547</td>
<td>5,078</td>
<td>45,172</td>
<td></td>
<td></td>
</tr>
<tr>
<td>512</td>
<td>1,765</td>
<td>2,047</td>
<td>3,734</td>
<td>18,078</td>
<td></td>
<td></td>
</tr>
</tbody>
</table>

TABLE 3 Run-times obtained by applying the Earley parser for context-free grammars to sentences in \( L_k \).
7.7 Conclusions

In this paper, we have presented a system that compiles parsing schemata to executable implementations of parsers, and used it to evaluate the performance of several TAG parsing algorithms, establishing comparisons both between themselves and with CFG parsers.

The results show that both string length and grammar size can be important factors in performance, and the interactions between them sometimes make their influence hard to quantify. The influence of string length in practical cases is usually below the theoretical worst-case bounds (between $O(n)$ and $O(n^2)$) in our tests for CFGs, and slightly below $O(n^3)$ for TAGs). Grammar size becomes the dominating factor in large TAGs, making tree filtering techniques advisable in order to achieve faster execution times.

Using TAGs to parse context-free languages causes an overhead both in constant factors and in practical computational complexity, thus increasing execution times by several orders of magnitude with respect to CFG parsing.

Acknowledgments

The work reported in this article has been supported in part by Ministerio de Educación y Ciencia and FEDER (TIN2004-07246-C03-01, TIN2004-07246-C03-02), Xunta de Galicia (PGIDIT05PXIC30501PN, PGIDIT05PXIC10501PN, Rede Galega de Procesamento da Linguaxe e Recuperación de Información), and Programa de becas FPU (Ministerio de Educación y Ciencia).

References


Properties of binary transitive closure logics over trees

STEPHAN KEPSER

Abstract

Binary transitive closure logic (FO∗ for short) is the extension of first-order predicate logic by a transitive closure operator of binary relations. Deterministic binary transitive closure logic (FO∗D) is the restriction of FO∗ to deterministic transitive closures. It is known that these logics are more powerful than FO on arbitrary structures and on finite ordered trees. It is also known that they are at most as powerful as monadic second-order logic (MSO) on arbitrary structures and on finite trees. We will study the expressive power of FO∗ and FO∗D on trees to show that several MSO properties can be expressed in FG∗ (and hence FO∗).

The following results will be shown.

• A linear order can be defined on the nodes of a tree.
• The class EVEN of trees with an even number of nodes can be defined.
• On arbitrary structures with a tree signature, the classes of trees and finite trees can be defined.
• There is a tree language definable in FO∗D that cannot be recognized by any tree walking automaton.
• FO∗ is strictly more powerful than tree walking automata.

These results imply that FO∗D and FO∗ are neither compact nor do they have the Löwenheim-Skolem-Upward property.

Keywords  Binary transitive closure logic

8.1 Introduction

The question about the best suited logic for describing tree properties or defining tree languages is an important one for model theoretic syntax as well as for querying treebanks. Model theoretic syntax is a research program in mathematical linguistics concerned with studying the descriptive complex-
ity of grammar formalisms for natural languages by defining their derivation trees in suitable logical formalisms. Since the very influential book by Rogers (1998) it is monadic second-order logic (MSO) or even more powerful logics that are used to describe linguistic structures.

With the advent of XML and query languages for XML documents, in particular XPath, the interest in logics for querying treebanks rose dramatically. There is now a large interest in this topic in computer science. Independent of that, but temporarily parallel, large syntactically annotated treebanks became available in linguistics. They provide nowadays a rich and important source for the study of language. But in order to access this source, suitable query languages for treebanks are required.

One of the simplest properties that are known to be inexpressible in first-order predicate logic (FO henceforth) is the transitive closure of a binary relation. It is therefore a natural move to extend FO by a binary transitive closure operator. And this move has been done before in the definition of query languages for relational databases, in particular for the SQL3 standard. But it seems that the expressive power of FO plus binary transitive closures (FO* for short) to define tree properties is not much studied yet. This is somewhat surprising, because there is reason to believe that FO* is more user friendly than MSO. Most users of query languages, in particular linguists, understand the concept of a transitive closure very well and know how to use it. It is a lot more difficult to use set variables to describe tree properties. An example for this claim is the fact MSO is capable of defining binary transitive closures, as shown by Moschovakis (1974). A formula expressing the transitive closure in MSO is given at the end of the next section. It is questionable that ordinary users (without profound knowledge of MSO) would be able to find this formula.

There exists a more restricted version of transitive closure, namely deterministic transitive closure (FO*D). The deterministic transitive closure of a binary relation is the transitive closure of the functional or deterministic part of the relation. We propose to seriously consider FO*D as a language for defining tree properties. We do so by showing that several important MSO definable properties can be defined in FO*D. One such example is the ability to define a linear order on the nodes of a tree. The order resembles depth-first left-to-right traversal of a tree. A linear order is a powerful concept that can be used defining additional properties. For example, it is used to count the number of nodes in a tree modulo a given natural number. An instance is the definition of the class EVEN of all trees with an even number of nodes in FO*D.

Arguably an important reason for Rogers’ choice of MSO is its ability to axiomatize trees. I.e., there exists a set of axioms such that an arbitrary structure (of a suitable signature) is a tree – finite or infinite – iff it is a model
of the axioms. It is known that this characterization of trees cannot be done using FO. But the full expressive power of MSO may not really be needed for the axiomatization, because we show that arbitrary trees and finite trees can be axiomatized in FO\(^*\). This capability of axiomatizing finite and infinite trees implies that FO\(^*\) (and hence also FO\(^*\)) is neither compact nor does it possess the Löwenheim-Skolem-Upward property.

There exists a tree automaton concept that defines serial instead of parallel processing of nodes in a tree, namely tree walking automata (TWA). As the name implies, a tree is processed by walking up and down in it and inspecting nodes serially. One may therefore believe that these automata could be the automaton-theoretic correspondent of FO\(^*\). But we show here that FO\(^*\) is more powerful. Every tree language that is recognized by a TWA can be defined in FO\(^*\). The relationship towards FO\(^D\)* is less clear. There are FO\(^D\)*-definable tree languages that cannot be recognized by any TWA.

### 8.2 Preliminaries

Let \( M \) be a set. We write \( \wp(M) \) for the power set of \( M \). Let \( R \subseteq M \times M \) be a binary relation over \( M \). The transitive closure \( TC(R) \) of \( R \) is the smallest set containing \( R \) and for all \( x, y, z \in M \) such that \( (x, y) \in TC(R) \) and \( (y, z) \in TC(R) \) we have \( (x, z) \in TC(R) \). I.e.,

\[
TC(R) := \bigcap \{ W \mid R \subseteq W \subseteq M \times M, \forall x, y, z \in M : (x, y), (y, z) \in W \implies (x, z) \in W \}.
\]

Deterministic transitive closure is the transitive closure of a deterministic, i.e., functional relation. For an arbitrary binary relation \( R \) we define its deterministic reduct by

\[
R_D := \{(x, y) \in R \mid \forall z : (x, z) \in R \implies y = z \}.
\]

Now

\[
DTC(R) := TC(R_D).
\]

We consider labeled ordered unranked trees. A tree is ordered if the set of child nodes of every node is linearly ordered. A tree is unranked if there is no relationship between the label of a node and the number of its children. For brevity we just write tree for labeled ordered unranked tree. In Sections 8.3 and 8.5 we only consider finite trees, in Section 8.4 we also consider infinite trees.

**Definition 1** A tree domain is a non-empty subset \( T \subseteq \mathbb{N}^* \) such that for all \( u, v \in \mathbb{N}^* \) : \( uv \in T \implies u \in T \) (closure under prefixes) and for all \( u \in \mathbb{N}^* \) and \( i \in \mathbb{N} \) : \( ui \in T \implies uj \in T \) for all \( j < i \) (closure under left sisters).

Let \( L \) be a set of labels. A tree is a pair \((T, Lab)\) where \( T \) is a tree domain and \( Lab : T \to L \) is a node labeling function.

A tree is finite iff its tree domain is finite.
We remark that a tree domain is at most countable, since it is a subset of a countable union of countable sets.

The languages to talk about trees will be extensions of first-order logic. Their syntaxes is as follows. Let \( L = \{ x, y, z, w, u, x_1, x_2, x_3, \ldots \} \) be a denumerable infinite set of variables. The atomic formulae are \( L(x) \) for each label \( L \in L \), \( x \rightarrow y \), \( x \downarrow y \), and \( x = y \). Complex formulae are constructed from simpler ones by means of the boolean connectives, existential and universal quantification, and transitive closure. I.e., if \( \phi \) and \( \psi \) are formulae, then \( \neg \phi, \phi \land \psi, \phi \lor \psi, \exists x \phi, \forall x \phi, \) and \( [\text{TC}_{x_1,x_2} \phi](x,y), [\text{DTC}_{x_1,x_2} \phi](x,y), \) respectively, are formulae.

The semantics of the first-order part of the language is standard. Let \( (T, \text{Lab}) \) be a tree. A variable assignment \( a : X \rightarrow T \) assigns variables to nodes in the tree. The root node has the empty address \( e \). Now \( \llbracket L(x) \rrbracket^a = T \iff \text{Lab}(a(x)) = L \), \( \llbracket x \downarrow y \rrbracket^a = T \iff a(y) = a(x) \) for some \( i \in \mathbb{N} \), i.e., \( \downarrow \) is the parent relation. \( \llbracket x \rightarrow y \rrbracket^a = T \iff \) there is a \( u \in T \) and \( i \in \mathbb{N} \) such that \( a(x) = ui \) and \( a(y) = ui + 1 \), i.e., \( \rightarrow \) is the immediate sister relation.

Boolean connectives and quantification have their standard interpretation. Now, \( \llbracket [\text{TC}_{x_1,x_2} \phi](x,y) \rrbracket^a = T \iff (a(x), a(y)) \in \text{TC}(\{(b,d) \mid \llbracket \phi \rrbracket^{ab/x_1d/x_2} = T\}) \)

where \( ab/x_1d/x_2 \) is the variable assignment that is identical to \( a \) except that \( x_1 \) is assigned to \( b \) and \( x_2 \) to \( d \). If \( \phi \) is a formula with free variables \( x_1, x_2 \), it can be regarded as a binary relation \( \phi(x_1, x_2) \). Then \( [\text{TC}_{x_1,x_2} \phi] \) is the transitive closure of this binary relation. This language is abbreviated \( \text{FO}^* \).

And \( \llbracket [\text{DTC}_{x_1,x_2} \phi](x,y) \rrbracket^a = T \iff (a(x), a(y)) \in \text{DTC}(\{(b,d) \mid \llbracket \phi \rrbracket^{ab/x_1d/x_2} = T\}) \).

This language is abbreviated \( \text{FO}^{**} \). It is simple to see that everything expressible in \( \text{FO}^{**} \) can also be expressed in \( \text{FO}^* \), because

\[
[\text{DTC}_{x_1,x_2} \phi(x_1, x_2)](x,y) \Leftrightarrow \ [\text{TC}_{x_1,x_2} \phi(x_1, x_2) \land \forall z \phi(x_1, z) \rightarrow z = x_2](x,y).
\]

It is an open question whether there are tree languages definable in \( \text{FO}^* \) that cannot be defined in \( \text{FO}^{**} \).

\( \text{FO}^* \) is amongst the smallest extension of first-order logic. It is known that the transitive closure of a binary relation is not first-order definable (Fagin, 1975). But when talking about trees, people frequently want to talk about paths in a tree. And a path is the transitive closure of certain base steps. \( \text{FO}^{**} \) and \( \text{FO}^* \) have at most the expressive power of monadic second-order logic (MSO). It is an old result, which goes back at least to Moschovakis (1974, p. 20), that the transitive closure of every MSO-definable binary relation is also MSO-definable. The following formula is due to Courcelle (1990). Let
$R$ be an MSO-definable binary relation. Then
\[
\forall X \left( \forall z, w \in X \land R(z, w) \implies w \in X \right) \land \forall z (R(x, z) \implies z \in X)
\]
is a formula with free variables $x$ and $y$ that defines the transitive closure of $R$.

It follows that every tree language definable in FO$^*$ can be defined in MSO. Whether the two logics are equivalent, seems an open question. For FO$^{po}$, the question is settled. Recently, Bojanczyk et al. (2006) have shown that the expressive power of MSO for defining tree languages properly extends the expressive power of FO$^{po}$.

### 8.3 Definability of Order

One of the abstract insights from descriptive complexity theory is that order is a very important property of structures. The relationship between certain logics and classical complexity classes is frequently restricted to ordered structures, i.e., structures where the carrier is linearly ordered. The reason for this restriction is to be found in the fact that computation is an ordered process. Definability and non-definability results for certain logics over ordered structures frequently do not extend to unordered structures. It is therefore an important property of a logic, if the logic itself is capable of expressing order without recourse to an extended signature. The probably best known logic with this property is $\Sigma^1_1$, the extension of first-order logic by arbitrary relation variables that are existentially quantified. It is obviously possible to define order in $\Sigma^1_1$ because we can say there is a binary relation that has all the properties of a linear order. These properties are known to be first-order properties. It is hence the ability to say “there is a binary relation” that is the key.

There is no way that FO$^{po}$ or FO$^*$ could define order on arbitrary finite structures. But if we only consider ordered trees as models, FO$^{po}$ can define order. Indeed it is possible to give a definition of the depth-first left-to-right order of nodes in a tree (and some variants).

**Proposition 9** There is an explicit definition of a linear order of the nodes in a tree in FO$^{po}$.

**Proof.** Define the proper dominance relation of trees $\text{Dom}(x, y)$ as $[\text{DTC}_{y,x} x \downarrow y](y, x)$. The idea of how to define dominance deterministically by walking upwards from the descendants to the ancestors goes back to Etessami and Immerman (1995). Similarly but simpler, define the sister relation $\text{Sis}(x, y)$ as $[\text{DTC}_{x,y} x \rightarrow y](x, y)$. Now define $x < y$ as

$$\text{Dom}(x, y) \lor (\exists w, v : \text{Sis}(w, v) \land (w = x \lor \text{Dom}(w, x)) \land (v = y \lor \text{Dom}(v, y)))$$

The first disjunct expresses the “depth-first” part of the order. The more complicated second disjunct formalizes the “left-to-right” part. It expresses that
there is a common ancestor of nodes \( x \) and \( y \) and node \( x \) is to be found on a left branch while \( y \) is to be found on a right branch. Care is taken that mutual domination is excluded. Hence the two disjuncts are mutually exclusive. Since the dominance and the sisterhood steps are both irreflexive, the whole relation \(<\) is irreflexive. Furthermore for each pair of distinct nodes in a tree, either one dominates the other, or there is a common ancestor such that one node is on a left branch while the other is on a right branch. Hence the relation is total. Transitivity can easily be checked by considering the four cases involved in expanding \( x < y \) and \( y < z \).

The proposition basically states that ordered trees are ordered structures in any logic at least as powerful as \( \text{FO}^\text{D} \). Note that the root node is the smallest element of the order. If the tree is finite, the largest element is the leaf of the rightmost branch of the tree. The root node is \( \text{FO} \)-definable via \( \neg \exists y : y \downarrow x \). The largest element \( \text{Max} \) of the order is \( \text{FO}^\text{D} \)-definable by \( \exists x \neg \exists y : x < y \). The successor \( y \) of a node \( x \) in the linear order \( (\text{Succ}(x, y)) \) is also \( \text{FO}^\text{D} \)-definable: \( x < y \land \neg \exists z : x < z \land z < y \). Using a linear order it is possible to count modulo some natural number on trees. That is for \( n, k \in \mathbb{N} \) we can define the class of finite trees such that each tree in the class has \( d \times n + k \) nodes (for some \( d \in \mathbb{N} \)). As an example, we define the class \( \text{EVEN} \) of trees with an even number of nodes (i.e, \( n = 2, k = 0 \)).

**Proposition 10** The class of finite trees with an even number of nodes is \( \text{FO}^\text{D} \)-definable.

**Proof.** We only consider the case where a tree has more than two nodes. The formula

\[
\exists w : \text{Succ}(\text{Root}, w) \land [\text{DTC}_{x,z} : \text{Succ}(x, z) \land \text{Succ}(z, y)](w, \text{Max})
\]

expresses that we go in one step from the root to its successor \( w \). From \( w \) we can reach the last element of the order by an arbitrary number of two successor steps. If we take the two-successors-step path through the linear order from the root to the maximum, we have an odd number of nodes, since a path of \( n \) double-successor-steps has \( n + 1 \) nodes. \( \square \)

**Corollary 11** \( \text{FO}^\text{D} \) has no normal form of the type \( [\text{DTC}_{x,z} \emptyset(x,y)](r,r) \) where \( \emptyset(x,y) \) is an FO formula and \( r \) the root. The same is true mutatis mutandis for \( \text{FO}^\text{D} \).

**Proof.** With a single application of a DTC-operator we can define trees with a linear order. If FO with a single DTC-operator is interpreted over finite successor structures, then it is equivalent to FO with order. But over finite orderings, \( \text{EVEN} \) is not definable in FO. \( \square \)

The above corollary is stated here because it contrasts with a fundamental result in descriptive complexity theory. Let \( \text{FO}(\text{TC}) \) be the extension of FO by
transitive closure operators of arbitrary width, that is the transitive closure of binary relations on tuples of arbitrary width. Let FO(DTC) be its deterministic counterpart. Immerman (1999) showed that both FO(TC) and FO(DTC) on ordered structures have a normal form consisting of a single outer application of the (deterministic) transitive closure operator on an otherwise FO formula.

8.4 Definability of Tree Structures

In previous and all following sections we assume that we only consider tree models as defined in the Preliminaries section. But in this section we take a more general view, a view that has its origin in model theoretic syntax. The aim is to find whether it is possible to give an axiomatization of those structures linguists are interested in. This task has two subparts. The first consists of defining trees, or more precisely finite trees, as the intended models. The second part consists of axiomatizing linguistic principles such as the Binding theory in the given logic. We will only be concerned with the first part here. This section is inspired by the book by Rogers (1998). More specifically we show that the main results of Chapter 3 carry over to $\text{FO}^\ast$. We will frequently cite this chapter in the current section.

The language of this section is deterministic binary transitive closure logic with equality over the following base relations:

- $\langle$ parent relation
- $\star$ dominance relation
- $+$ proper dominance relation
- $<$ left-of relation

We also assume there to be a set $\mathcal{L}$ of unary predicate symbols representing linguistic labels. We write $\text{FO}^\ast$ for this language to indicate that the base relations differ from the ones in the other sections of this paper.

A model for $\text{FO}^\ast$ is a tuple $(U, P, D, L, \text{Lab})$ where $U$ is a non-empty domain, $P, D$ and $L$ are binary relations over $U$ interpreting $\langle, \star$ and $\ast$. And $\text{Lab} : \mathcal{L} \to \wp(U)$ interprets each label as a subset of $U$.

Since the intended models of this language are trees, we have to restrict the class of models by giving axioms of trees. Many properties of trees can be defined by first-order axioms. The following 12 axioms are cited from (Rogers, 1998, p. 15f.).

A1 $\exists x \forall y : x \ast y$
  (Connectivity . dominance)
A2 $\forall x, y : (x \ast y \land y \ast x) \rightarrow x = y$
  (Antisymmetry of dominance)
A3 $\forall x, y, z : (x \ast y \land y \ast z) \rightarrow x \ast z$
  (Transitivity of dominance)
A4 $\forall x, y : x \ast y \leftrightarrow (x \ast y \land x \neq y)$


A5 \( \forall x, y : x \triangleright y \leftrightarrow (x \triangleright x' y \land \forall z : (x \triangleright x' z \land z \triangleright y) \rightarrow (z \triangleright x' x \lor x' y)) \)

(Definition of proper dominance)

A6 \( \forall x, z : x \triangleright x' z \rightarrow ((\exists y : x \triangleright y \land y \triangleright z) \land (\exists y : y \triangleright z)) \)

(Discreteness of dominance)

A7 \( \forall x, y : (x \triangleright x' y \land y \triangleright x) \leftrightarrow (x \not\triangleright y \land y \not\triangleright x) \)

(Exhaustiveness and exclusiveness)

A8 \( \forall w, x, y : (x \triangleright x' w \land y \triangleright y' z) \rightarrow w \prec z \)

(Inheritance of Left-of wrt. dominance)

A10 \( \forall w, x, y : (x \prec y \land y \prec z) \rightarrow x \prec z \)

(Transitivity of left-of)

A11 \( \forall x : (\exists y : x \triangleright y) \rightarrow (\exists z : x \triangleright y \land z \triangleright z' y) \)

(Existence of a minimum child)

A12 \( \forall x, z : x \prec z \rightarrow (\exists w : x \prec w \land \forall y : w \prec z \Rightarrow y \not\prec w) \land (\exists z : y \prec z \land \forall w : w \prec z \Rightarrow y \not\prec w) \)

(Discreteness of left-of)

A discussion of these axioms can be found in (Rogers, 1998, p. 16f.). Every tree (finite or infinite) obeys to these axioms. But there are non-standard models, i.e., structures that are models of these axioms but would not be considered as trees. Actually, it is not possible to give a first-order axiomatization of trees, as was shown by Backofen et al. (1995). The simplest example of a non-standard model can be gained by adapting the well-known example of a non-standard model of FO arithmetics to tree structures. This model is depicted in Figure 1. The carrier is the disjoint union of the natural numbers and the integers. The dominance relation is defined by taking the natural order on natural numbers and integers plus every natural number dominates every integer. Formally: 

\[
U = \mathbb{N} \cup \mathbb{Z}, \quad P = \{(n, n+1) \mid n \in \mathbb{N} \cup \mathbb{Z}\}, \quad D = \{(n, m) \mid n, m \in \mathbb{N}, n \leq m \} \cup \{(n, m) \mid n, m \in \mathbb{Z}, n \leq m \} \cup \{(n, z) \mid n \in \mathbb{N}, z \in \mathbb{Z}\}, \quad L = \emptyset.
\]

This model is not a tree because the integers are infinitely far away from the root.

The FO axioms demand that the proper dominance relation does not only contain the immediate dominance relation but also the transitive closure of the immediate dominance. In the non-standard model, proper dominance truly extends the transitive closure of immediate dominance. All natural numbers properly dominate all integers. But this part of the dominance relation is not contained in the transitive closure of immediate dominance. In a proper tree model, the proper dominance is always identical to the transitive closure of immediate dominance. This insight can be expressed in FO\(^\omega_w\) as an axiom.

AT1 \( \forall x, y : x \triangleright y \rightarrow [\text{DTC}_{w, z} w](y, x) \)

(Proper dominance is the transitive closure of immediate dominance)
Another way of reading this axiom is to say that the path from an arbitrary node back to the root is finite.

AT1 together with the first-order axioms does still not suffice to axiomatize proper trees. An example of a non-standard model for which AT1 holds true is given in Figure 2. Formally, we set \( U = \{ r \} \cup \mathbb{N} \cup \mathbb{Z} \), \( P = \{ (r, z) \mid z \in \mathbb{N} \cup \mathbb{Z} \} \), \( D = P \cup \{ (i, i) \mid i \in \{ r \} \cup \mathbb{N} \cup \mathbb{Z} \} \), and \( L = \{ (n, m) \mid n, m \in \mathbb{N}, n < m \} \cup \{ (n, m) \mid n, m \in \mathbb{Z}, n < m \} \cup \{ (n, z) \mid n \in \mathbb{N}, z \in \mathbb{Z} \} \). Consider the sisters of a node. They are ordered by \( \prec \), and there is a left-most sister. Now, in a proper tree, the number of sisters to the left is finite for every node. In the model in Figure 2 all integers have infinitely many left sisters. This configuration has to be avoided by means of one more axiom as follows. We can easily define that one node is the immediate sister of another node. The relation \( IS(x, y) \) is defined as \( \exists z : z \prec x \land z \prec y \land x < y \land \neg \exists w : x < w < y \). Now we can spell out an axiom analog to AT1.
Theorem 12 Axioms A1–A12, AT1, and AT2 define the class of tree models.

The proof is analogous to the proof of Theorem 3.9 in (Rogers, 1998). Consider in particular Footnote 8 on page 23.

Proof. Rogers showed that every tree (in the sense of Definition 1) is a model of axioms A1–A12 and for each node \( x \in U \) the sets \( A_x = \{(y, x) \in D\} \) of ancestors of \( x \) and \( L_x = \{y \mid \exists z : (z, x), (z, y) \in D \text{ and } (y, x) \in L\} \) of left sisters of \( x \) are finite (Lemma 3.5). And every tree obviously satisfies axioms AT1 and AT2.

Furthermore, each model of axioms A1–A12 where \( A_x \) and \( L_x \) are finite for each node \( x \in U \) is isomorphic to a tree (Lemma 3.6).

Now suppose a model of A1–A12 satisfies AT1. Then for each node \( x \in U \) the set \( A_x \) is finite, because it contains the root (A1) and is constructed of parent-child steps (AT1), and a transitive closure of single steps cannot reach a limit ordinal. An analogous argument can be made with respect to models of A1–A12 and AT2. Hence for every model of of A1–A12, AT1, and AT2 and all nodes \( x \in U \) we see that the sets \( A_x \) and \( L_x \) are finite. By the above quoted Lemma 3.6, these models are isomorphic to trees. \(\square\)

The tree models of Axioms A1–A2, AT1, and AT2 can be finite as well as infinite. But since they are all tree models, they are at most countable. This is because every tree domain is at most countable (see remark after Definition 1). And every tree model is isomorphic to a tree. As an immediate consequence we get that FO\(^*\) does not have the Löwenheim-Skolem-Upward property. This property states that if a theory (i.e., potentially infinite set of sentences) has a model of size \( \omega \) it has models of arbitrary infinite cardinalities. It is a typical property of FO logic.

Corollary 13 The logics FO\(^*\) and FO\(^+\) do not have the Löwenheim-Skolem-Upward property.

Linguists are mostly (if not exclusively) concerned with finite trees. Hence it would be nice if we could restrict the class of models further down to finite trees. This can indeed be done. Rogers (1998) defines a linear order on the nodes of a tree as follows. Node \( x < y \) iff \( x \not< y \vee x < y \). By Axiom A7, each pair of nodes is either a member of the dominance relation or a member of the left-of relation. Hence this defines indeed a linear order. Actually, the order is the same as the one in the previous section: depth-first left-to-right tree traversal. As in the previous section we use \( Succ(x, y) \) for \( y \) being the immediate successor of \( x \) in the order. Finiteness can now be defined in two steps. Firstly we demand the linear order to be the deterministic transitive closure of the immediate successor relation. The consequence of this demand is that
for every element in the order there is only a finite number of nodes that are smaller than this element. Secondly we demand the order to have a maximal element. If the maximal element has only a finite number of elements smaller than it, the tree is obviously finite.

\[ \forall x, y : x < y \implies [\text{DTC}_{<} x, y] \wedge \exists x \forall y : y < x \vee y = x. \]  
(Finiteness of the order \(<\))

**Theorem 14** Axioms A1–A12, AT1, AT2, and AF define the class of finite tree models.

**Proof.** By Theorem 12, every model of the Axioms A1–A12, AT1, and AT2 is isomorphic to a tree model. If a model is finite, then AF is certainly true. For the converse, assume that \( \forall x, y : x < y \implies [\text{DTC}_{<} x, y] \). By definition of the DTC-operator, the set \( \{ y \mid y < x \} \) of elements smaller than \( x \) is finite for every node \( x \). If the order has additionally a maximal element \( m \), then it is finite. \( \square \)

This theorem implies that another property of FO, namely compactness, does not extend to FO\(^{D^*}\).

**Corollary 15** The logics FO\(^{D^*}\) and FO\(^{*}\) are not compact.

FO, on the other hand, is not capable of defining the class of finite trees. It is well known that compactness and definability of finiteness of models mutually exclude each other.

### 8.5 Transitive Closure Logics and Tree Walking Automata

Tree walking automata were introduced by Aho and Ullman (1971) as sequential automata on trees. At every moment of its run, a TWA is in a single node of the tree and in one of a finite number of states. It walks around the tree choosing a neighboring node based on the current state, the label of the current node, and the child number of the current node.

More formally, we consider trees of maximal branching degree \( k \). The following definition is mainly cited from (Bojanczyk and Colcombet, 2005). Every node \( v \) has a type. The possible values are \( \text{Types} = \{ r, 1, 2, \ldots, k \} \times \{ l, i \} \) where \( r \) stands for the root, \( j \in \{1, \ldots, k\} \) states that \( v \) is the \( j \)-th child, \( l \) states that \( v \) is a leaf, \( i \) that \( v \) is an internal node. A direction is an element of \( \text{Dir} = \{ \uparrow, \downarrow_1, \ldots, \downarrow_k, \text{stay} \} \) where \( \uparrow \) stands for ‘move to the parent’, \( \downarrow_j \) ‘move to the \( j \)-th child, and \( \text{stay} \) to ‘stay at the current node’. A TWA is a quintuple \( (S, \Sigma, \delta, s_0, F) \) where \( S \) is a finite set of states, \( \Sigma \) is the alphabet of node labels, \( s_0 \in S \) is the initial state and \( F \subseteq S \) is the set of final states. The transition relation \( \delta \) is of the form

\[ \delta \subseteq (S \times \text{Types} \times \Sigma) \times (S \times \text{Dir}). \]
A configuration is a pair of a node and a state. A run is a sequence of configurations where every two consecutive configurations are consistent with the transition relation. A run is accepting if it starts and ends at the root of the tree, the first state is $s_0$ and the last state is a member of $F$. The TWA accepts a tree if there is an accepting run. The set of $\Sigma$-trees recognized by a TWA is the set of trees for which there is an accepting run.

Bojanczyk and Colcombet (2005) showed that TWA cannot recognize all regular tree languages. This means that MSO and tree automata are strictly more powerful than TWA. In an extension of their proof we will show that even FO$^*$ is more powerful than TWA.

**Theorem 16** The classes of tree languages definable in FO$^*$ strictly extend the classes of tree languages recognizable by TWA.

**Proof.** The proof consists of two parts. We will first show that every TWA-recognizable tree language is FO$^*$-definable. Secondly we will show that there is an FO$^*$-definable tree language that cannot be recognized by any TWA.

The first part of the proof is based on recent results by Neven and Schwentick (2003). They showed that a tree language is recognizable by a TWA if and only if it is definable by a formula of the following type: $\text{[TC}_{x,y} \phi(x,y))(r,r)$ where $r$ is a constant for the root of a tree, $\phi$ is an FO formula with additional unary $\text{depth}_m$ predicates. Apart from the $\text{depth}_m$ predicates, these formulae are obviously in FO$^*$. Now, $\text{depth}_m(x)$ is true iff $x$ is a multiple of $m$ steps away from the root. For every $m$, the predicate $\text{depth}_m$ can be defined by an FO$^*$-formula: $[\text{TC}_{t_0,t_m} \exists x_1,\ldots,x_{m-1} : x_0 \downarrow x_1 \land \cdots \land x_{m-1} \downarrow x_m](r,x)$ is a predicate that is true on a node $x$ just in case there is a $k \in \mathbb{N}$ such that $x$ is at depth $k \times m$. Thus every TWA-recognizable tree language is FO$^*$-definable.

To show the second half of the theorem, we will indicate that the separating language $L$ given by Bojanczyk and Colcombet (2005) can be defined in FO$^*$. The authors consider binary trees. They show (in Fact 1) that $L$ can be defined in first-order logic with the following three basic relations: left and right child, and ancestor relation. Now, left and right child are obviously FO$^*$-definable relations. And the ancestor relation is – as in the previous sections – FO$^*$-definable by $[\text{DTC}_{y,x} x \downarrow y](y,x)$. $\Box$

**Corollary 17** There exists an FO$^*$-definable tree language that is not TWA-recognizable.

Please note that there exists an alternative proof of Theorem 16. Engelfriet and Hoogeboom (2006) have recently shown that transitive closure logics correspond to certain pebble automata. (A pebble automaton is a TWA enhanced by a finite sets of pebbles.) More precisely, the deterministic pebble automata have exactly the same expressive power as deterministic binary transitive clo-
Properties of binary transitive closure logics over trees / 115

Sure logic. And non-deterministic pebble automata have the same expressive power as binary transitive closure logic where each transitive closure operator is under the scope of an even number of negations. Since a TWA is a pebble automaton with 0 pebbles, the first half of above theorem follows from the equivalence results of (Engelfriet and Hoogeboom, 2006). The second half of the theorem follows from new results by Bojanczyk et al. (2006) who show that each additional pebble extends the expressive power of a pebble automaton. Bojanczyk et al. (2006) also provide an alternative proof of Corollary 17. As a result, either TWA and DPA are incomparable, or TWA are less powerful than DPA.

8.6 Conclusion

We showed a range of properties of FO∗∗ and FO∗ to indicate that they should seriously be considered as logics for defining tree languages. Although the addition of binary transitive closure to first-order logic can be seen as a small one, FO∗∗ is capable of expressing important second-order properties over trees. It is possible to define a linear order over the nodes in a tree. And using this order one can count modulo any natural number. On arbitrary structures with appropriate signature one can axiomatize the classes of trees and finite trees. These axiomatizations showed that FO∗∗ is neither compact nor does it have the Löwenheim-Skolem-Upward property. Furthermore although tree walking automata look like they might serve as an automaton model for FO∗, it turns out that FO∗ is more powerful than TWA.

A word about complexity issues may be in place. FO∗∗ and FO∗ have quite a good data complexity. By translating FO∗ formulae into MSO formulae and using the equivalence between MSO and tree automata one can see that FO∗ has a linear time data complexity. And since FO∗ is a sub-logic of FO(TC), it also has NLOGSPACE data complexity whereas FO∗∗ has LOGSPACE data complexity. A straight-forward implementation of transitive closure yields a PTIME query complexity. It is unclear to the author whether this result can be improved upon.

The main open question is of course whether FO∗ is strictly less powerful than MSO. It is also interesting to study the relationship of FO∗ to modal languages for trees like PDL_{Tree} (Kracht, 1995). Marx (2004) basically showed that PDL_{Tree} is at most as powerful as FO∗∗, where FO∗∗ is the restriction of FO∗ where every formula has at most 3 different variables. ten Cate (2006) recently showed that queries in XPath with Kleene star and loop predicate have the same expressive power as FO∗.

One may also ask what happens if we introduce the transitive closure of arbitrary relations, not just binary ones. This logic (abbreviated FO(TC)) was introduced by Immerman (see Immerman, 1999) to logically describe
the complexity class NLOGSPACE. Tiede and Kepser (2006) have recently shown that FO(TC) is more expressive than MSO over trees. The statement remains true even if one only considers deterministic transitive closures.

Acknowledgments
The author wishes to thank four anonymous referees whose comments helped improving the quality of the paper. This research was funded by a grant of the German Research Foundation (DFG SFB-441).

Stephan Kepser
Collaborative Research Centre 441
University of Tübingen
Germany

References


Pregroups with modalities
ALEKSANDRA KISLAK-MALINOWSKA

Abstract
In this paper we concentrate mainly on the notion of $\beta$-pregroups, which are pregroups (first introduced by Lambek (1999) in 1999) enriched with modality operators. $\beta$-pregroups were first proposed by Fadda (2002) in 2001. The motivation to introduce them was to (locally) limit the associativity in the calculus considered. In this paper we present this new calculus in the form of a rewriting system, and prove the very important feature of this system: that in a given derivation the non-expanding rules must always precede non-contracting ones in order for the derivation to be minimal (normalization theorem). We also propose a sequent system for this calculus and prove the cut elimination theorem for it.

Keywords Pregroup, $\beta$-pregroup, normalization theorem, cut elimination

9.1 Introduction

Definition 2 A pregroup is a structure $(G, \leq, \cdot, l, r, 1)$ such that $(G, \leq, \cdot, 1)$ is a partially ordered monoid, and $l, r$ are unary operations on $G$, fulfilling the following conditions:
\[ a'l \leq 1 \leq a'a \leq 1 \leq a'a \] (9.1)
for all $a \in G$. Element $a'$ ($a'$ respectively) is called the left (right) adjoint of $a$.

The notion of a pregroup, introduced by Lambek (1999), is connected to the notion of a residuated monoid, known from the theory of partially ordered algebraic systems.

Theorem 18 (Lambek (1999)) In each pregroup the following equalities and inequalities are valid:
\[ 1' = 1 = 1', \quad a'r = a = ar, \] (9.2)
\[(ab)^{i} = b^{i}a^{i}, \quad (ab)^{j} = b^{j}a^{j}, \quad (9.3)\]

\[a \leq b \iff b^{i} \leq a^{i} \iff b^{j} \leq a^{j}. \quad (9.4)\]

For any arbitrary element \(a\) of a pregroup we define an element \(a^{(n)}\), for \(n \in \mathbb{Z}\), in the following way:

\[a_{0} = a, \quad a^{(n+1)} = (a^{(n)})^{r}, \quad a^{(n-1)} = (a^{(n)})^{l}. \quad (9.4)\]

As a consequence of (2) and (9.4) we obtain:

\[a^{(n)}a^{(n+1)} \leq 1 \leq a^{(n+1)}a^{(n)} \quad (9.5)\]

\[\text{if } a \leq b \text{ then } a^{(2n)} \leq b^{(2n)} \text{ and } b^{(2n+1)} \leq a^{(2n+1)} \quad (9.6)\]

for all \(n \in \mathbb{Z}\).

Let \((P, \leq)\) be a poset. Elements of the set \(P\) are treated as constants. Terms are expressions of the form \(p^{(n)}\), for \(p \in P, n \in \mathbb{Z}\); \(p^{(0)}\) is equal to \(p\). Types are finite strings of terms, denoted by \(X, Y, Z, V, U\) etc. The basic rewriting rules are as follows:

- **(CON)** - contraction:
  \[X, p^{(n)}, p^{(n+1)}, Y \rightarrow X, Y;\]

- **(EXP)** - expansion:
  \[X, Y \rightarrow X, p^{(n+1)}, p^{(n)}, Y;\]

- **(IND)** - induced step:
  \[X, p^{(2n)}, Y \rightarrow X, q^{(2n)}, Y, \quad X, q^{(2n+1)}, Y \rightarrow X, p^{(2n+1)}, Y, \quad \text{for } p \leq q \text{ in } (P, \leq).\]

Furthermore, we consider derivations \(X \Rightarrow Y\) in \(F(P)\) (free pregroup generated by \((P, \leq)\)). Following Lambek (2001), we distinguish two special cases:

- **(GCON)** - generalized contraction:
  \[X, p^{(2n)}, q^{(2n+1)}, Y \rightarrow X, Y; \quad X, q^{(2n-1)}, p^{(2n)}, Y \rightarrow X, Y, \quad \text{where } p \leq q \text{ in } (P, \leq).\]

- **(GEXP)** - generalized expansion:
  \[X, Y \rightarrow X, q^{(2n+1)}, q^{(2n)}, Y; \quad X, Y \rightarrow X, q^{(2n)}, p^{(2n-1)}, Y, \quad \text{where } p \leq q \text{ in } (P, \leq).\]

The relation \(\Rightarrow\) is a reflexive and transitive closure of the relation \(\rightarrow\).

**Theorem 19 (Lambek switching lemma, Lambek (1999))** If \(X \Rightarrow Y\) is in \(F(P)\), then there exist types \(U, V\) such that we can go from type \(X\) to \(U\) \((X \Rightarrow U)\) using only generalized contractions, from type \(U\) to \(V\) \((U \Rightarrow V)\) using only induced steps, and from type \(V\) to \(Y\) \((V \Rightarrow Y)\) using only generalized expansions.

From the above mentioned lemma we obtain:

**Corollary 20 (Buszkowski (2003))** If \(X \Rightarrow Y\) in \(F(P)\), and \(Y\) is a simple type or an empty string, then \(X\) can be transformed into \(Y\) only by means of (CON).
and (IND). If $X \Rightarrow Y$ in $F(P)$, and $X$ is a simple type or an empty string, then $X$ can be transformed into $Y$ only by means of (EXP) and (IND).

## 9.2 Pregroups with modalities

In this section we generalize some definitions and results concerning pregroups introduced in Lambek (1999). The definition of a pregroup with $\beta$-operator was proposed by Fadda (2002). The motivation to introduce modality operators stems from the fact there was a need to (locally) limit associativity in the calculus considered.

**Definition 3** A pregroup with $\beta$-operator is a pregroup $G$ enriched additionally with a monotone mapping $\beta : G \rightarrow G$.

**Definition 4** A $\beta$-pregroup is a pregroup with $\beta$-operator such that $\beta$-operator has the right adjoint $\hat{\beta}$ ($\hat{\beta}$-operator), i.e., there exists a monotone mapping $\hat{\beta} : P \rightarrow P$ with the property that for all $a$ and $b$ in $P$, $\beta(a) \leq b$ if and only if $a \leq \hat{\beta}(b)$.

It is easy to show that $\hat{\beta}$-operators, if they exist, are uniquely defined and connected to $\beta$-operators with the following rules of expansion and contraction, for all $a \in P$.

\[ a \leq \hat{\beta}(\beta(a)) \quad \text{and} \quad \beta(\hat{\beta}(a)) \leq a. \]  

(9.7)

The basic rewriting rules are as follows:

1. **Contracting rules**
   - $(CON)$ - contraction:
     \[ X, p^n, p^{(n+1)}, Y \rightarrow X, Y; \]
   - $(B - CON)$ - B-contraction:
     \[ X, [B(Y)]^{(\alpha)}_{(\alpha+1)}, [B(Y)]^{(\alpha+1)}_{(\alpha+1)}, Z \rightarrow X, Z; \]
     where $B \in \{\beta, \hat{\beta}\}$.
   - $(\beta - CON)$ - $\beta$-contraction:
     \[ X, [\beta(\hat{\beta}(Y))]^{(n)}_{(n+1)}, Z \rightarrow X, Y^{(n+1)}_{(n+1)}, Z; \]
     \[ X, [\hat{\beta}(\beta(Y))]^{(n+1)}_{(n+1)}, Z \rightarrow X, Y^{(n)}_{(n+1)}, Z; \]
   - $(B - IND_c)$ - $B_c$ induced step:
     \[ X, [B(Y)]^{(n)}_{(n+1)}, Z \rightarrow X, [B(Y)]^{(n+1)}_{(n+1)}, Z; \]
     where $B \in \{\beta, \hat{\beta}\}$, and $Y_1 \rightarrow Y_2$ is a contracting rule.
     \[ X, [B(Y)]^{(n+1)}_{(n+1)}, Z \rightarrow X, [B(Y)]^{(n+1)}_{(n+1)}, Z; \]
     where $B \in \{\beta, \hat{\beta}\}$, a $Y_1 \rightarrow Y_2$ is an expanding rule.

2. **Expanding rules**
   - $(EXP)$ - expansion:
     \[ X, Y \rightarrow X, p^n_{(n+1)}, p^n_{(n+1)}, Y; \]
   - $(B - EXP)$ - B-expansion:
     \[ X, Z \rightarrow X, [B(Y)]^{(n)}_{(n+1)}, [B(Y)]^{(n+1)}_{(n+1)}, Z; \]
     where $B \in \{\beta, \hat{\beta}\}$. 
\[ (\beta - \text{EXP}) - \beta - \text{expansion:} \]
\[ X, Y^{(2n)}, Z \rightarrow X, [\beta(\beta(Y))]^{(2n)}, Z; \]
\[ X, Y^{(2n+1)}, Z \rightarrow X, [\beta(\beta(Y))]^{(2n+1)}, Z. \]

\[ (B - \text{IND}_\beta) - B_\beta \text{ induced step:} \]
\[ X, [B(Y_1)]^{(2n)}, Z \rightarrow X, [B(Y_2)]^{(2n)}, Z; \]
where \( B \in [\beta, \beta] \), a \( Y_1 \rightarrow Y_2 \) is an expanding rule.
\[ X, [B(Y_2)]^{(2n+1)}, Z \rightarrow X, [B(Y_1)]^{(2n+1)}, Z; \]
where \( B \in [\beta, \beta] \), a \( Y_1 \rightarrow Y_2 \) is a contracting rule.

3. **P-rules (neither expanding nor contracting)**

\[ (\text{IND}) - \text{induced step:} \]
\[ X, p^{(2n)}, Y \rightarrow X, q^{(2n)}, Y, \]
\[ X, q^{(2n+1)}, Y \rightarrow X, p^{(2n+1)}, Y, \quad \text{for } p \leq q \text{ w} (P, \leq). \]

\[ (B - \text{IND}_P) - B_P \text{ induced step:} \]
\[ X, [B(Y_1)]^{(2n)}, Z \rightarrow X, [B(Y_2)]^{(2n)}, Z; \]
where \( B \in [\beta, \beta] \), and \( Y_1 \rightarrow Y_2 \) is a P-rule.
\[ X, [B(Y_2)]^{(2n+1)}, Z \rightarrow X, [B(Y_1)]^{(2n+1)}, Z; \]
where \( B \in [\beta, \beta] \), and \( Y_1 \rightarrow Y_2 \) is a P-rule.

In the above mentioned rules we assume that \( p, q \) are elements of \( P \), whereas \( X, Y, Z, Y_1, Y_2 \) are elements of \( P' \). The relation \( \Rightarrow \) is a reflexive and transitive closure of the relation \( \rightarrow \).

Fadda (2002) gives some examples illustrating the usage of \( \beta \) - pregroups for natural language. Among others, he shows that assigning a type \([\beta(X)]X[\beta(X)] \) to the conjunction and (where \( X \) is an arbitrary type), will let us see the structure of a sentence more clearly.

Consider the sentence: John and Mary left. Applying the calculus of pregroups without modalities we can show that the string of types assigned to given words can be reduced to the type of a sentence. However, the order of consecutive contraction is important here (\( np \) means a noun phrase):

\[ (*) \quad \text{John and Mary left.} \]
\[ np \quad np' \quad np' \quad np' \quad np' \quad np' s \quad \rightarrow \]
\[ np' \quad np' \quad np' s \quad \rightarrow \]
\[ np \quad np' s \quad \rightarrow \quad s \]

\[ (**) \quad \text{John and Mary left.} \]
\[ np \quad np' \quad np' \quad np' \quad np' \quad np' s \quad \rightarrow \]
\[ np' \quad np' \quad np' s \quad \rightarrow \]
\[ np' \quad np' s \quad \rightarrow \quad s \]

In the second case \( (**) \) we do not get a type \( s \). Applying the calculus of \( \beta \)-pregroups, we could handle the above mentioned sentence in the following way:

\[ (*** \quad \text{John and Mary left.} \]
\[ \beta(np) \quad [\beta(np)]p np [\beta(np)]^l \quad \beta(np) \quad np' s \quad \rightarrow \quad s \]
In that case the structure of types `induces’ the order of contractions.

**Normalization theorem for β - pregroups**

Further we consider derivations of a type $X \Rightarrow Y$.

**Definition 5** A derivation is called non-expanding, if there are no expanding rules present.

**Definition 6** A derivation is called non-contracting, if there are no contracting rules present.

**Definition 7** Composition of derivations $d_1(X \Rightarrow U)$ and $d_2(U \Rightarrow Y)$ is a derivation $Y$ from $X$, which transforms first $X$ into $U$ according to $d_1$, and then $U$ into $Y$ according to $d_2$.

**Definition 8** A derivation $d(X \Rightarrow Y)$ is called normal, if it is a composition of some non-expanding derivation $d_1(X \Rightarrow U)$ and some non-contracting derivation $d_2(U \Rightarrow Y)$.

On elements of $P'$ we introduce a measure in the following way:

$$
\mu(x) = 0,
\mu(p^{(p)}) = 1,
\mu(B(Y)) = \mu(Y) + 1, \text{ for } B \in \{\beta, \hat{\beta}\}
\mu(Y_1, ..., Y_k) = \mu(Y_1) + ... + \mu(Y_k).
$$

A measure on the rewriting rules is defined as follows:

$$
\mu(CON) = 2,
\mu(EXP) = 2,
\mu(\beta - CON) = 2,
\mu(\beta - EXP) = 2,
\mu(B - CON) = 2 + 2\mu(Y),
\mu(B - EXP) = 2 + 2\mu(Y),
\mu(IND) = 1,
\mu(B_c - IND) = 1 + \mu(d(Y_1 \rightarrow Y_2)),
\mu(B_e - IND) = 1 + \mu(d(Y_1 \rightarrow Y_2)),
\mu(B_p - IND) = 1 + \mu(d(Y_1 \rightarrow Y_2)),
\mu(d(X_0 \Rightarrow X_i)) = \mu(d(X_0 \rightarrow X_1)) + ... + \mu(d(X_{i-1} \rightarrow X_i)),
$$

where $X_0 \Rightarrow X_i$ means $X_0 \rightarrow X_1 \rightarrow ... \rightarrow X_i$.

**Definition 9** A derivation $d(X \Rightarrow Y)$ is called minimal, if it has the least possible measure of all derivations $Y$ from $X$, and the least possible complexity (which is understood as a sum of measures of all rules used in the derivation).

**Definition 10** The position of a given rule in the derivation $X_0 \rightarrow X_1 \rightarrow ... \rightarrow X_i$ is number $i$, such that $X_{i-1} \rightarrow X_i$ is the occurrence of this rule in the derivation.
Definition 11  A degree of non-normal derivation \(d(X \Rightarrow Y)\) is the minimal position of a contracting rule which occurs (not necessarily directly) after an expanding rule. A degree of normal derivation is number 0.

Theorem 21 (Normalization theorem for \(\beta\)-pregroups)  Every minimal derivation is normal.

Proof. Let \(X_0 \rightarrow X_1 \rightarrow \ldots \rightarrow X_0\) be a minimal derivation. Let \(i\) be a degree of this derivation. We will show that \(i = 0\), and as a consequence our derivation is normal. Assume that \(i > 0\). Of course \(1 < i \leq n\) from the definition of a degree. Let \(j\) be the greatest number less than \(i\), such that \(X_{j-1} \Rightarrow X_j\) is the occurrence of an expanding rule.

Let \(R_1\) denote the rule used on the position \(j\), and \(R_2\) the rule used in the position \(i\). The following cases are to be considered:

1. \(R_1 = (\text{EXP})\) \(R_2 = (\text{CON})\),
2. \(R_1 = (\text{EXP})\) \(R_2 = (B - \text{CON})\),
3. \(R_1 = (\text{EXP})\) \(R_2 = (\beta - \text{CON})\),
4. \(R_1 = (\text{EXP})\) \(R_2 = (B - \text{IND}_0)\),
5. \(R_1 = (B - \text{EXP})\) \(R_2 = (\text{CON})\),
6. \(R_1 = (B - \text{EXP})\) \(R_2 = (B - \text{CON})\),
7. \(R_1 = (B - \text{EXP})\) \(R_2 = (\beta - \text{CON})\),
8. \(R_1 = (B - \text{EXP})\) \(R_2 = (B - \text{IND}_0)\),
9. \(R_1 = (\beta - \text{EXP})\) \(R_2 = (\text{CON})\),
10. \(R_1 = (\beta - \text{EXP})\) \(R_2 = (B - \text{CON})\),
11. \(R_1 = (\beta - \text{EXP})\) \(R_2 = (\beta - \text{CON})\),
12. \(R_1 = (\beta - \text{EXP})\) \(R_2 = (B - \text{IND}_0)\),
13. \(R_1 = (B - \text{IND}_0)\) \(R_2 = (\text{CON})\),
14. \(R_1 = (B - \text{IND}_0)\) \(R_2 = (B - \text{CON})\),
15. \(R_1 = (B - \text{IND}_0)\) \(R_2 = (\beta - \text{CON})\),
16. \(R_1 = (B - \text{IND}_0)\) \(R_2 = (B - \text{IND}_0)\).

In the proof of this theorem the above mentioned cases are considered. In all cases we assume that the rule \(R_1\) occurs on the position \(j\), and the rule \(R_2\) on the position \(i\). All steps \(X_j \rightarrow X_{j+1} \rightarrow \ldots \rightarrow X_{i-1}\) consist of application of non-expanding and non-contracting rules. These must be of the form of either \((\text{IND})\) or \((B_p - \text{IND})\). None of these steps cannot be independent from \(X_{j-1} \Rightarrow X_j\), as otherwise we could do the last of independent steps after \(R_2\), getting the derivation with the same measure but the lower degree. We can also assume that none of these steps is not independent from \(X_{j-1} \Rightarrow X_j\); otherwise it would transform our derivation performing the first step before \(R_1\), increasing the number \(j\), and changing neither \(i\) nor \(\mu(d(X \Rightarrow Y))\).

If the rules \(R_1\) and \(R_2\) are adjacent (without intermediate P-rules), we change the order in case they are independent from each other (getting the derivation of smaller complexity); in case they are dependent from each other,
we show that this part of derivation can be transformed using rules of smaller complexity - thus showing that the initial derivation was not normal.

Considering the sixteen cases mentioned above, we show that non-expanding rules must always precede non-contracting ones. Otherwise our derivation would not be minimal, which would be a contradiction to our assumption. Thus every minimal derivation must be normal.

As the proof is long and technical, we show as an example only one of above mentioned sixteen cases:

Case 1.1. \( R_1 = (\text{EXP}) \quad R_2 = (\text{CON}) \),

\[ X_{i-1} \rightarrow X_i \] is of the form \( S, T \rightarrow S, p^{(n+1)}, p^{(n)}, T \); \( X_{i-1} \rightarrow X_i \) is of the form \( U, q^{(n)}, q^{(n+1)}, V \rightarrow U, V \). The derivation \( X_{i-1} \rightarrow X_i \rightarrow \ldots \rightarrow X_{i-1} \rightarrow X_i \) could be as follows:

\[ S, p_0^{(2n)}, T \rightarrow S, p_0^{(2n)}, p_k^{(2n)}, T \rightarrow S, p_0^{(2n)}, p_k^{(2n)}, p_k^{(2n)}, T \rightarrow \ldots \rightarrow S, p_0^{(2n)}, p_k^{(2n)}, T \rightarrow S, p_k^{(2n)}, T \], (assuming \( p_0 \leq p_1 \leq \ldots \leq p_k \)). The measure of a new derivation is \( \mu(d(X_{i-1} \Rightarrow X_i)) = k \) (k times the rule (IND) was used). We reach a contradiction, as the measure of the second derivation is smaller. We showed that the initial derivation was not normal.

\[ \square \]

**Corollary 22** If \( X \Rightarrow Y \) in a free \( \beta \)-pregroup, and \( Y \) is a simple type or an empty string, then \( Y \) can be derived from \( X \) only by means of non-expanding rules.

If \( X \Rightarrow Y \) in a free \( \beta \)-pregroup, and \( X \) is a simple type or an empty string, then \( Y \) can be derived from \( X \) only by means of non-contracting rules.

### 9.3 Axiom system for pregroups with modalities

The rewriting system given in the previous section can also be presented as the calculus of sequents in a Gentzen style. Let \( (P, \leq) \) be fixed. Atoms and types are defined as before. Sequents are of the form \( X \Rightarrow Y \), where \( X, Y \) are types. The axiom and inference rules are as follows:

<table>
<thead>
<tr>
<th>Rule</th>
<th>Premises</th>
<th>Conclusion</th>
</tr>
</thead>
<tbody>
<tr>
<td>(Id)</td>
<td>( X \Rightarrow X )</td>
<td>( X,Y \Rightarrow Z )</td>
</tr>
<tr>
<td>(RA)</td>
<td>( X \Rightarrow Y, Z )</td>
<td>( X,Y \Rightarrow Z )</td>
</tr>
<tr>
<td>(LIND)</td>
<td>( X, p^{(n)}, q^{(n+1)}, Y \Rightarrow Z )</td>
<td>( X, p^{(n)}, Y \Rightarrow Z )</td>
</tr>
<tr>
<td>(RIND)</td>
<td>( X \Rightarrow Y, p^{(n)}, Z )</td>
<td>( X \Rightarrow Y, p^{(n)}, Z )</td>
</tr>
<tr>
<td>( X, p^{(n+1)}, Y \Rightarrow Z )</td>
<td>( X \Rightarrow Y, q^{(n)}, Z )</td>
<td></td>
</tr>
<tr>
<td>( X, q^{(n+1)}, Y \Rightarrow Z )</td>
<td>( X \Rightarrow Y, q^{(n+1)}, Z )</td>
<td></td>
</tr>
<tr>
<td>( X, q^{(n+1)} \Rightarrow Z )</td>
<td>( X \Rightarrow Y, q^{(n+1)}, Z )</td>
<td></td>
</tr>
</tbody>
</table>

In rules (LIND) and (RIND) we assume that \( p \leq q \) in \( P \). \( X, Y, Z \) are any
arbitrary types, \( p, q \) are arbitrary elements of \( P \), for \( n \in \mathbb{Z} \).

(\text{BLA}) \begin{align*}
X, T & \Rightarrow Z \\
X, [B(Y)]^{(n)}, [B(Y)]^{(n+1)}, T & \Rightarrow Z
\end{align*}

(\text{BRA}) \begin{align*}
X & \Rightarrow T, Z \\
X, B(Y)[(n+1)], T & \Rightarrow Z
\end{align*}

\[ \frac{X, Y^{n+1}, T \Rightarrow Z}{X \Rightarrow T, Y^{2n+1}, Z} \quad (\beta \text{ LA}) \]

\[ \frac{X, [B(Y)]^{(2n)}, T \Rightarrow Z}{X \Rightarrow T, [B(Y)]^{(2n)}, Z} \quad (\beta \text{ RA}) \]

In rules (BLA), (BRA), (BLIND) and (BRIND), \( B \in \{\beta, \beta\} \). Additionally, in rules (BLIND) we assume that \( Y \rightarrow Y_2 \) arises as a result of a non-expanding rule in an even case, and a non-contracting rule in an odd case, in a rewriting system from a former section. In rules (BRIND) we assume that \( Y \rightarrow Y_2 \) arises as a result of non-contracting rule in an even case, and non-expanding rule in an odd case, in a rewriting system form a former section.

The cut rule is of the form

(\text{CUT}) \begin{align*}
X & \Rightarrow Y, Y \Rightarrow Z \\
X & \Rightarrow Z
\end{align*}

Let \( MS \) denote the system axiomatized by (Id), (LA), (RA), (LIND), (RIND), (BLA), (BRA), (\( \beta \)-LA), (\( \beta \)-RA), (BLIND) and (BRIND). Let \( MS' \) denote the system \( MS \) enriched additionally with a cut rule (CUT).

### 9.3.1 Cut elimination for the systems with modalities

We show that for above mentioned systems the following theorems hold:

**Theorem 23** For all types \( X, Y \), \( X \Rightarrow Y \) holds in the sense of a rewriting system if and only if \( X \Rightarrow Y \) is provable in \( MS' \).

**Proof.** Assume \( X \Rightarrow Y \) holds in the sense of the rewriting system. Then, there exist types \( Z_0, \ldots, Z_n \), \( n \geq 0 \), such that \( Z_0 = X, Z_n = Y \), and \( Z_{i-1} \rightarrow Z_i \), \( 1 \leq i \leq n \). We show that \( Z_{i-1} \Rightarrow Z_i \) is provable in \( MS' \), for \( 1 \leq i \leq n \). (Here we show it only for a few chosen cases.)

1. If \( Z_{i-1} \rightarrow Z_i \) is the case of (CON), so it is of the form \( X, p^{(n)}, p^{(n+1)}, Y \rightarrow X, Y \), we apply (LA) to axiom \( X, Y \Rightarrow X, Y \). We get
   \[ X, Y \Rightarrow X, Y \]
   \[ X, p^{(n)}, p^{(n+1)}, Y \Rightarrow X, Y \]

7. If \( Z_{i-1} \rightarrow Z_i \) is the case of (IND), so it is of the form:

7.1. \( X, p^{(2n)}, Y \rightarrow X, q^{(2n)}, Y \), for \( p \leq q \), we apply (LIND) to axiom

\[ X, p^{(2n)}, Y \Rightarrow X, q^{(2n)}, Y \]

We get

\[ X, q^{(2n)}, Y \Rightarrow X, q^{(2n)}, Y \]

We can also apply (RIND) to axiom \( X, p^{(2n)}, Y \Rightarrow X, p^{(2n)}, Y \). We obtain

\[ X, p^{(2n)}, Y \Rightarrow X, p^{(2n)}, Y \]

\[ X, q^{(2n)}, Y \Rightarrow X, q^{(2n)}, Y \]

\[ X, q^{(2n)}, Y \Rightarrow X, q^{(2n)}, Y \]

\[ X, q^{(2n)}, Y \Rightarrow X, q^{(2n)}, Y \]
Theorem 24 (Cut elimination theorem) For all types $X, Y$, $X \Rightarrow Y$ is provable in $MS'$ if and only if $X \Rightarrow Y$ is provable in $MS$.

Proof. The ‘only if’ part is obvious. If for all types $X, Y$, $X \Rightarrow Y$ is provable in $MS$ (without CUT), it is also provable in $MS'$.

Assume that $X \Rightarrow Y$ is provable in $MS'$. By the theorem 23, $X \Rightarrow Y$ holds in the rewriting system. From the theorem 21 there exists such type $U$, that $X \Rightarrow U$ holds only by using non-expanding rules, whereas $U \Rightarrow Y$ holds only by using non-contracting rules. Thus, there exist types $Z_0, \ldots, Z_m$, ($m \geq 0$), such that $Z_0 = X, Z_m = U$ and for all $1 \leq i \leq m, Z_{i-1} \Rightarrow Z_i$ is a result of non-expanding rules. We show that $Z_i \Rightarrow U$ is provable in $MS$, for all $0 \leq i \leq m$. $Z_m \Rightarrow U$ is an axiom (Id). Assume that $Z_i \Rightarrow U$ is provable in $MS, i > 0$. If $Z_{i-1} \rightarrow Z_i$ is (CON), then $Z_{i-1} \Rightarrow U$ is a result of applying (LA) to $Z_i \Rightarrow U$. If $Z_{i-1} \rightarrow Z_i$ is $(B \Rightarrow \text{CON})$, then $Z_{i-1} \Rightarrow U$ is a result of applying ($\beta$LA) to $Z_i \Rightarrow U$. If $Z_{i-1} \rightarrow Z_i$ is (IND), then $Z_{i-1} \Rightarrow U$ is a result of application (LIND) to $Z_i \Rightarrow U$. If $Z_{i-1} \rightarrow Z_i$ is $(B \Rightarrow \text{IND})$, then $Z_{i-1} \Rightarrow U$ is a result of applying (BLIND) to $Z_i \Rightarrow U$. If $Z_{i-1} \rightarrow Z_i$ is $(B \Rightarrow \text{IND}_p)$, then $Z_{i-1} \Rightarrow U$ is a result of applying (BLIND) to $Z_i \Rightarrow U$.

Now, there exist types $V_0, \ldots, V_n$, $n \geq 0$, such that $V_0 = U, V_n = Y$, an for all $1 \leq i \leq n, V_{i-1} \rightarrow V_i$ is a result of applying a non-contracting rule. We show that $X \Rightarrow V_i$ is provable in $MS$, for all $0 \leq i \leq n$. $X \Rightarrow V_0$ is provable in $MS$ from the first part of the proof. Assume that $X \Rightarrow V_{i-1}$ is provable in
MS, $1 \leq i$. If $V_{i-1} \rightarrow V_i$ is $(EXP)$, then $X \Rightarrow V_i$ is a result of applying (RA) to $X \Rightarrow V_{i-1}$. If $V_{i-1} \rightarrow V_i$ is $(B \rightarrow EXP)$, then $X \Rightarrow V_i$ is a result of applying (BRA) to $X \Rightarrow V_{i-1}$. If $V_{i-1} \rightarrow V_i$ is $(\beta \rightarrow EXP)$, then $X \Rightarrow V_i$ is a result of applying ($\beta$RA) to $X \Rightarrow V_{i-1}$. If $V_{i-1} \rightarrow V_i$ is $(IND)$, then $X \Rightarrow V_i$ is a result of applying (RIND) to $X \Rightarrow V_{i-1}$. If $V_{i-1} \rightarrow V_i$ is $(B \rightarrow IND)$, then $X \Rightarrow V_i$ is a result of applying (BRIND) to $X \Rightarrow V_{i-1}$.

Thus, we showed that $X \Rightarrow Y$ is provable in MS. 

\[ \square \]

9.4 Conclusion

In this paper we presented pregroups with modalities. First, we presented them in the form of a rewriting system, then we proposed the sequent system for them and finally showed the connections between those two presentations. Using those connections we were able to prove the cut elimination theorem.

References


Simpler TAG semantics through synchronization

REBECCA NESSON AND STUART SHIEBER

Abstract

In recent years Laura Kallmeyer, Maribel Romero, and their collaborators have led research on TAG semantics through a series of papers refining a system of TAG semantics computation. Kallmeyer and Romero bring together the lessons of these attempts with a set of desirable properties that such a system should have. First, computation of the semantics of a sentence should rely only on the relationships expressed in the TAG derivation tree. Second, the generated semantics should compactly represent all valid interpretations of the input sentence, in particular with respect to quantifier scope. Third, the formalism should not, if possible, increase the expressivity of the TAG formalism. We revive the proposal of using synchronous TAG (STAG) to simultaneously generate syntactic and semantic representations for an input sentence. Although STAG meets the three requirements above, no serious attempt had previously been made to determine whether it can model the semantic constructions that have proved difficult for other approaches. In this paper we begin exploration of this question by proposing STAG analyses of many of the hard cases that have spurred the research in this area. We reframe the TAG semantics problem in the context of the STAG formalism and in the process present a simple, intuitive base for further exploration of TAG semantics. We provide analyses that demonstrate how STAG can handle quantifier scope, long-distance WH-movement, interaction of raising verbs and adverbs, attitude verbs and quantifiers, relative clauses, and quantifiers within prepositional phrases.

Keywords SYNCHRONOUS TREE-ADJOINING GRAMMAR, STAG SEMANTICS

10.1 Introduction

In recent years Laura Kallmeyer, Maribel Romero, and their collaborators have led research on TAG semantics through a series of papers refining a system of TAG semantics computation using evolving techniques including enriched derivation tree structure (Kallmeyer, 2002a,b), flexible composition
of feature-based TAG with a semantic representation associated with each elementary tree (Kallmeyer and Joshi, 2003, Joshi et al., 2003, Kallmeyer, 2003), semantic features in a more expressive extension of feature-based TAG (Gardent and Kallmeyer, 2003), and, most recently, semantic features on the derivation tree itself (Kallmeyer and Romero, 2004, Romero et al., 2004). Kallmeyer and Romero (2004) bring together the lessons of these attempts with a set of desirable properties that such a system should have. First, computation of the semantics of a sentence should rely only on the relationships expressed in the TAG derivation tree. Because TAG elementary trees represent minimal semantic units, the only information necessary for semantic computation should be the information encoded in the derivation tree: which elementary trees have combined and the address at which the combining operation took place. Second, the generated semantics should compactly represent all valid interpretations of the input sentence, in particular with respect to quantifier scope. Third, the formalism should not, if possible, increase the expressivity of the TAG formalism.

We revive the proposal of using synchronous TAG (STAG) to simultaneously generate syntactic and semantic representations for an input sentence (Shieber and Schabes, 1990). Although STAG meets the three requirements above, no serious attempt had previously been made to determine whether it can model the semantic constructions that have proved difficult for other approaches. In this paper we begin exploration of this question by proposing STAG analyses of many of the hard cases that have spurred the research in this area. We reframe the TAG semantics problem in the context of the STAG formalism and in the process present a simple, intuitive base for further exploration of TAG semantics.

After reviewing STAG in Section 10.2, we provide analyses in Sections 10.3.1 through 10.3.4 for sentences that exemplify several hard cases for TAG semantics that have been raised by Kallmeyer and others in recent papers: quantifier scope (as exemplified by sentences (17) and (21), presented below along with the desired semantic interpretations), long-distance WH-movement (18), interaction of raising verbs and adverbs, attitude verbs and quantifiers (19,20,21), relative clauses (22), and quantifiers within prepositional phrases (23) (Kallmeyer and Romero, 2004, Romero et al., 2004, Joshi et al., 2003, Kallmeyer, 2003, Kallmeyer and Joshi, 2003).\(^1\)

\begin{equation}
\text{(17) Everyone likes someone.}
\end{equation}

\begin{align}
\text{every}(x, \text{person}(x), \text{some}(z, \text{person}(z), \text{like}(x, z))) \\
\text{some}(z, \text{person}(z), \text{every}(x, \text{person}(x), \text{like}(x, z)))
\end{align}

\begin{equation}
\text{(18) Who does Bill think Paul said John likes?}
\end{equation}

\begin{equation}
\text{who}(y, \text{think}(\text{bill}, \text{say}(\text{paul}, \text{like}(\text{john}, y))))
\end{equation}

\(^1\)We notate curried two-place relations \(P(x)(y)\) as \(P(y, x)\) for readability.
Simpler TAG semantics through synchronization

(19) Bill thinks John apparently likes Mary.
\[\text{think}(\text{bill}, \text{apparently}(\text{like}(\text{john}, \text{mary})))\]

(20) John sometimes likes everyone.
\[\text{every}(x, \text{person}(x), \text{sometimes}(\text{like}(\text{john}, x)))\]
\[\text{sometimes}(\text{every}(x, \text{person}(x), \text{like}(\text{john}, x)))\]

(21) Bill thinks everyone likes someone.
\[\text{think}(\text{bill}, \text{every}(x, \text{person}(x), \text{some}(z, \text{person}(z), \text{like}(x, z))))\]
\[\text{think}(\text{bill}, \text{some}(z, \text{person}(z), \text{every}(x, \text{person}(x), \text{like}(x, z))))\]

(22) A problem whose solution is difficult stumped Bill.
\[\text{a}(x, \text{and}(\text{problem}(x)), \text{the}(y, \text{and}(\text{solution}(y), \text{poss}(x, y)), \text{isDifficult}(y))), \text{stumped}(\text{bill}, x))\]

(23) Two politicians spy on someone from every city.
\[\text{two}(x, \text{politician}(x), \text{every}(z, \text{city}(z), \text{some}(y, \text{person}(y) \land \text{from}(z, y)), \text{spyOn}(x, y))))\]
\[\text{every}(z, \text{city}(z), \text{some}(y, \text{person}(y) \land \text{from}(z, y)), \text{two}(x, \text{politician}(x), \text{spyOn}(x, y))))\]
\[\text{two}(x, \text{politician}(x), \text{some}(y, \text{every}(z, \text{city}(z), \text{person}(y) \land \text{from}(z, y)), \text{spyOn}(x, y))))\]
\[\text{some}(y, \text{every}(z, \text{city}(z), \text{person}(y) \land \text{from}(z, y)), \text{two}(x, \text{politician}(x), \text{spyOn}(x, y))))\]

10.2 Introduction to Synchronous TAG

A tree-adjoining grammar (TAG) consists of a set of elementary tree structures and two operations, substitution and adjunction, used to combine these structures. The elementary trees can be of arbitrary depth. Each internal node is labeled with a nonterminal symbol. Frontier nodes may be labeled with either terminal symbols or nonterminal symbols and one of the diacritics ↓ or *. Use of the diacritic ↓ on a frontier node indicates that it is a substitution node. The substitution operation occurs when an elementary tree rooted in the nonterminal symbol \(A\) is substituted for a substitution node labeled with the nonterminal symbol \(A\). Auxiliary trees are elementary trees in which the root and a frontier node, called the foot node and distinguished by the diacritic *, are labeled with the same nonterminal. The adjunction operation involves splicing an auxiliary tree with root and designated foot node labeled with a nonterminal \(A\) at a node in an elementary tree also labeled with nonterminal
A. Examples of the substitution and adjunction operations on sample elementary trees are shown in Figure 1.

Synchronous TAG (STAG) extends TAG by taking the elementary structures to be pairs of TAG trees with links between particular nodes in those trees. An STAG is a set of triples, \( \langle t_L, t_R, \sim \rangle \) where \( t_L \) and \( t_R \) are elementary TAG trees and \( \sim \) is a linking relation between nodes in \( t_L \) and nodes in \( t_R \) (Shieber, 1994, Shieber and Schabes, 1990). Derivation proceeds as in TAG except that all operations must be paired. That is, a tree can only be substi-

---

**FIGURE 1** Example TAG substitution and adjunction operations.

**FIGURE 2** An English syntax/semantics STAG fragment (a), derived tree pair (b), and derivation tree (c) for the sentence “John apparently likes Mary.”
tuted or adjoined at a node if its pair is simultaneously substituted or adjoined at a linked node. We notate the links by using boxed indices \[\text{marking linked nodes.}\]

Figure 2 contains a sample English syntax/semantics grammar fragment that can be used to parse the sentence “John apparently likes Mary”. The node labels we use in the semantics correspond to the semantic types of the phrases they dominate.\(^2\) Variables such as \(x\) in the semantic tree in Figure 3 are taken to be bound in the obvious way, so that in multiple uses of the tree they can be presumed to be renamed apart.

Figure 2(c) shows the derivation tree for the sentence. Substitutions are notated with a solid line and adjunctions are notated with a dashed line. Note that each link in the derivation tree specifies a link number in the elementary tree pair. The links provide the location of the operations in the syntax tree and in the semantics tree. These operations must occur at linked nodes in the target elementary tree pair. In this case, the noun phrases \(\text{John}\) and \(\text{Mary}\) substitute into \(\text{likes}\) at links \([\text{3}]\) and \([\text{4}]\) respectively. The word \(\text{apparently}\) adjoins at link \([\text{2}]\). The resulting semantic representation can be read off the derived tree by treating the leftmost child of a node as a functor and its siblings as its arguments. Our sample sentence thus results in the semantic representation \(\text{apparently}(\text{likes}(\text{john}, \text{mary}))\).

10.3 STAG Analyses of the Phenomena

10.3.1 Quantifier Scope and Wh-Words

For sentence (17), we would like to generate a scope-neutral semantic representation that allows both the reading where \(\text{some}\) takes scope over \(\text{every}\) and the reading where \(\text{every}\) takes scope over \(\text{some}\). We propose a solution in which a derivation tree with multiple adjunction nondeterministically determines multiple derived trees each manifesting explicit scope (Schabes and Shieber, 1993); the derivation tree itself is therefore the scope neutral representation.

The multi-component quantifier approach followed by Joshi et al. (2003) suggests a natural implementation of quantifiers in STAG.\(^3\) In this approach the syntactic tree for quantifiers has two parts, one that corresponds to the scope of the quantifier and attaches at the point where the quantifier takes scope, and the other that contains the quantifier itself and its restriction and attaches where syntactically expected at a noun phrase. In their work, a single-

\(^2\)This representation is for the sake of readability. The labels could be replaced using any well-chosen finite set of nonterminal symbols.

\(^3\)The multi-component approach to quantifiers in STAG was first suggested by Shieber and Schabes (1990) under the rewriting definition of STAG derivation where the order of rewriting produced the scope ambiguity. Williford (1993) explored the use of multiple adjunction to achieve scope ambiguity.
node auxiliary tree is used for the scope part of the syntax in order to get the desired relationship between the quantifier and the quantified expression in features threaded through the derivation tree and hence in the semantics. Using STAG, we do not need the single-node auxiliary tree in the syntax because we can pair the usual syntactic representation for quantified NPs with a multi-component semantic representation that expresses the same idea (Figure 3). In order to use these quantifiers, we change the links in the elementary trees for verbs to allow a single link to indicate two positions in the semantics where a tree pair can adjoin, as shown in Figure 3.\(^4\)

Given this representation of quantifiers we get the derivation tree shown in Figure 3 for sentence (17).\(^5\) Note that the resulting derivation tree neces-

\(^4\)We have chosen here to add the three-way links in addition to the existing links in the tree for unquantified noun phrases such as proper nouns (though we suppress the two-way NP links in the figures for readability). Another possibility would be to remove the two-way links. In this case, all noun phrases would be “lifted” à la Montague. That is, even unquantified noun phrases would have a scope part, which could be a single-node auxiliary tree.

\(^5\)We notate multi-component insertions that involve both a substitution and an adjunction with a combination dashed and dotted line.
sarily incorporates multiple adjunction (Schabes and Shieber, 1993), that is, multiple auxiliary trees are adjoined at the same node in an auxiliary tree. In particular, the scope parts of both every and some attach at the root of the semantic tree of likes. Such cases of multiple adjunction induce ambiguity; the derivation tree represents multiple derived trees. In the case at hand, the derivation is ambiguous as to which quantifier scopes higher than the other. This ambiguity in the derivation tree thus models the set of valid scopings for the sentence. In essence, this method uses multiple adjunction to model scope-neutrality.

This same method can be used to obtain the correct scope relations for sentences with long-distance WH-movement such as sentence (18) using the multi-component elementary tree pair for who and the elementary tree pairs for thinks (the tree pair for says is similar) and likes in the WH context given in Figure 4. Kallmeyer and Romero (2004) highlight this case as difficult because in the usual syntactic analysis there is no link in the derivation tree between who and thinks or between thinks and likes, but in the desired semantics who takes scope over the thinks proposition and the likes proposition is an argument to thinks.

In our analysis, by contrast, the semantics follows quite naturally from the standard syntactic analysis of the structure of the likes elementary tree in the WH context and the elementary tree pair for thinks given in Figure 4. The derivation of this sentence is also given in Figure 4. Note that it is required by the structure of the trees that who take scope over thinks.
FIGURE 5 Derivation trees for (a) “Bill thinks John apparently likes Mary”, (b) “John sometimes likes everyone”, and (c) “Bill thinks everyone likes someone.”

10.3.2 The Interaction Between Attitude Verbs, Raising Verbs, Adverbs and Quantifiers

The interaction between attitude verbs and raising verbs or adverbs as in sentences (19), (20), and (21) has been problematic for TAG semantics (Kallmeyer and Romero, 2004). A successful analysis must be flexible enough to produce the correct semantics for sentence (19) even though there is no link between thinks and apparently in the derivation tree. It must also be flexible enough to allow all scope orderings between VP modifiers and quantifiers as in sentence (20). In fact, given the elementary trees we have already presented and the ones for attitude verbs demonstrated by Figure 4, our analysis already allows for scope interactions among all these elements. Indeed, because the semantic components of attitude verbs, VP modifiers, and quantifiers all adjoin at the same node in the semantic tree of the verb, our analysis allows all scope orderings among them. This is clearly too permissive, because it allows quantifiers to scope out of the finite clause in which they appear and would allow a reading of sentence (19) in which apparently scopes over thinks. To prevent quantifiers from scoping out of the finite clause in which they appear, as in sentences (19) and (21), we can add an additional adjunction site to the semantic trees for verbs above the current root node. This is shown in Figure 6 in the likes tree pair. The link configuration ensures that attitude verbs (adjoining at link 1) will now scope higher than all VP modifiers (adjoining at 2) and quantifiers (adjoining at links 3 and 4). VP modifiers and quantifiers will still be able to take all scope orderings relative to each other. Using the modified verb trees, STAG produces the correct semantics for sentences (19), (20), and (21) with the derivations given in Figure 5.

10.3.3 Relative Clauses

Relative clauses provide another putatively difficult case for TAG semantics because both the main verb and the relative clause need access to the variable introduced by the determiner as in sentence (22) (Kallmeyer, 2003). We overcome this difficulty and compute the desired semantics by introducing higher-order functions into the semantic trees using lambda-calculus notation. This modification allows us to maintain tree-locality. The syntac-
tic analysis we use is similar to that of Kallmeyer (2003) in that it maintains the Condition on Elementary Tree Minimality (Frank, 1992) and uses the relative pronoun to introduce the relative clause. However, it treats the relative pronoun as a noun modifier rather than a noun phrase modifier. We also posit the existence of “lifted” versions of the elementary trees for verbs in which their argument positions have been abstracted over. We use a higher-order conjunction and that relates two properties: $\lambda PQx.P(x) \land Q(x)$, and a higher-order se function that relates two properties and makes use of the higher-order conjunction: $\lambda PQx.y.(P, \lambda z.\text{poss}(x, z))(y), Q(y))$. The elementary tree pairs and resulting derivation tree for sentence (22) are given in Figure 7. The derived tree is given in Figure 8. When reduced, the resulting semantics is $\lambda z.\lambda x.(\text{problem}(x) \land \text{the}(y, \text{solution}(y) \land \text{poss}(x, y), \text{isDifficult}(y))), \text{stumped}(\text{bill}, z))$. 

FIGURE 6 Modified tree for $\text{likes}$ that enforces a restriction on quantifiers scoping outside of the finite clause.

FIGURE 7 Key elementary trees and derivation for “A problem whose solution is difficult stumped Bill.”
10.3.4 Nested Quantifiers and Inverse Linking

Quantifiers in prepositional phrases such as in sentence (23) pose another challenge for TAG semantics (Joshi et al., 2003). Although a nested quantifier may take scope over the quantifier within which it is nested (so-called “inverse linking”) not all permutations of scope orderings of the quantifiers are available (Joshi et al., 2003). In particular, readings in which a quantifier intervenes between a nesting quantifier and its nested quantifier are not valid. In our example sentence (23), this predicts that the readings some > two > every and every > two > some should not be valid. Joshi et al. (2003) introduce a special device allowing nesting and nested quantifiers to form an indivisible quantifier set during the derivation, which prevents other quantifiers from intervening between them. In our solution, because the nested quantifier is introduced through the prepositional phrase, which in turn modifies the noun phrase containing the nesting quantifier, the two quantifiers already naturally form a set that operates as a unit with respect to the rest of the derivation.\(^6\) The elementary tree pairs and derivation trees for our analysis of (23) are shown in Figure 9.

One notable feature of this analysis is that the four different scope readings that result are not the product of a single derivation tree. The alternate scope orderings for the nested and nesting quantifier exist because there are two available adjunction sites for the scope of quantifiers in the prepositional phrase.
phrase to attach. This results in two distinct derivation trees. The alternate scope orderings for this quantifier set and the remaining quantifier are obtained by multiple adjunction at the root of the verb tree. The set of valid derivation trees for a sentence thus constitutes the scope neutral representation. This set of trees may be compactly represented, for instance as a shared forest.  

10.4 Comparison to the Kallmeyer and Romero Approach  

As mentioned above, research on TAG semantics has been led by Laura Kallmeyer, Maribel Romero, and their collaborators through a series of papers refining a system of TAG semantics computation using feature unification and other formal devices (Kallmeyer and Romero, 2004, Romero et al., 2004, Kallmeyer, 2003, Kallmeyer and Joshi, 2003, Joshi et al., 2003, Gardent and Kallmeyer, 2003). Although their approach has evolved over time, the underlying principles of using the relationships expressed in the derivation tree as the basis for the computation and generating underspecified semantic representations have been constant. In its current formulation, they perform semantic computation by attaching semantic feature structures directly to the nodes in the derivation tree. When carefully chosen, these features unify to produce an underspecified representation of the semantics of a sentence that, when further disambiguated, generates the set of valid interpretations. In one or another of their recent papers they have provided successful analyses of each of the hard cases that we have addressed here, though some of their analyses might have to be restated to bring them up to date with the newest formulation of their method.

This analysis, like that of Joshi et al. (2003), makes several predictions about quantifier scope that might be disputed. First, some argue that more than four scope orderings should be available for sentences like sentence (23) (VanLehn, 1978, Hobbs and Shieber, 1987). This analysis cannot generate additional scope orderings without breaking tree set locality. Second, the scope readings in which the nesting quantifier takes scope over the nested quantifier result in the nested quantifier having scope over the restriction of the nesting quantifier but not over its scope. Donkey sentence constructions such as “Every man with two books loves them” call this prediction into question.
Our work owes much to theirs both for the clear formulation of the problems and the progress in formulating analyses for some of the hard cases. The primary advantage of our approach is its conceptual simplicity. The clear separation of syntax and semantics, the directness of the link interface, and the familiarity of the TAG operations used in our approach make it very simple. The semantic-feature-unification-based approach has become cleaner and easier to understand as Kallmeyer and others have refined it over the years. Nonetheless, it is safe to say that the amount of formal machinery—including propositional labels, separate individual and propositional variables, semantic representations consisting of a set of formulas and a set of scope constraints, features on the derived tree and the derivation tree, each semantic feature structure containing a nested feature structure for each address in the elementary syntax tree, each of these feature structures containing features to handle binding of propositional and individual variables, feature unification, flexible composition, and quantifier sets—necessary to solve the range of problems that we have addressed here, is qualitatively more complex. In fact, we use no formal machinery that had not been introduced by 1994 in the TAG literature.

An additional advantage of our approach is that it does not increase the expressivity of the TAG formalism. One might think that the inclusion of multiple adjunction would lead to an increase in expressivity (Dras, 1999). However, because links can only be used once in an STAG derivation, only a finite number of multiple adjunctions may occur at a single adjunction site. This rules out problematic uses of multiple adjunction. Kallmeyer and Romero maintain the semantic features on the derivation tree rather than in the feature structures already used in the feature-based TAGs (FTAG) of their syntax in part because the set of semantic feature structures is not finite, potentially increasing the expressivity of the FTAG formalism (Kallmeyer and Romero, 2004). Although moving the features to the derivation tree avoids increasing the expressivity of the formalism used for syntax when taken alone, the additional expressivity in the features of the semantics could be used to block operations in the syntax thereby filtering the syntax to produce non-tree-adjointing languages. It remains to be seen whether this additional expressivity will be required for TAG semantics.

Advantages and disadvantages of the different methods aside, in this still nascent area of research it is desirable to have several quite different approaches at our disposal as we explore the hard problems presented by generating natural language semantics in the TAG framework. Our approach revives an old idea with the aim of opening a new avenue for research into semantics in the TAG framework.
10.5 Conclusion

We have presented the synchronous TAG formalism as a method for computing semantics in the TAG framework, and have shown that it enables simple, natural analyses for all of the cases that have exercised recent attempts at formulating formal semantics for TAG. It satisfies each of the desiderata laid out at the beginning of this paper. First, it does not require any additional information other than that available in the derivation tree to generate the semantics. Because the syntax and semantic representations are built up synchronously, the derivation tree set is a complete specification of the relationship between them. Nothing other than the set of elementary tree pairs and the synchronous TAG operations are required to generate a semantic representation. Second, the derivation tree set provides a compact representation for all valid semantic interpretations of the given sentence. Using multiply-adjoined quantifiers we take advantage of the ambiguity in the interpretation of the derivation tree that is introduced by multiple adjunction. We take each possible ordering of multiply-adjoined trees to be valid. We leave open the possibility of using an additional method to prefer certain scope orders and disprefer or eliminate others. Third, the STAG system, as used, does not increase the expressivity of the TAG formalism (Shieber, 1994). Finally, our analysis is a straightforward expression of a simple idea: we use TAG for both syntax and semantics and use the derivation tree and the links between trees in elementary tree pairs as the interface between them.

10.6 Acknowledgments

This work was supported in part by grant IIS-0329089 from the National Science Foundation. We wish to thank Rani Nelken and the three anonymous reviewers for valuable comments on earlier drafts.

References


Encoding second order string ACG with deterministic tree walking transducers

SYLVAIN SALVATI

Abstract
In this paper we study the class of string languages represented by second order Abstract Categorial Grammar. We prove that this class is the same as the class of output languages of deterministic tree walking automata. Together with the result of de Groote and Pogodalla (2004) this shows that the higher-order operations involved in the definition of second order ACGs can always be represented by operations that are at most fourth order.

Keywords Abstract categorial grammar, λ-calculus, deterministic tree walking transducers, mildly context sensitive languages

11.1 Introduction
Abstract Categorial Grammars (ACGs) (de Groote (2001)) are based on the linear logic (Girard (1987)) and on the linear λ-calculus. They describe the surface structures by using for syntax the ideas Montague (1974) devoted to semantics. ACGs describe parse structures with higher-order linear λ-terms and syntax as a higher-order linear homomorphism (lexicon) on parse structures. Intuitively, the higher the order of the parse structures is, the richer should the languages of analysis be and the higher the order of the lexicons is, the richer should the class of languages be. On the one hand, de Groote and Pogodalla (2004) have shown how to encode of several context free formalisms by using second order parse structures (i.e. sets of trees). They have encoded Context Free Grammars using second order lexicons, Linear Context Free Tree Grammars using third order lexicons and Linear Context Free Rewriting Systems (Weir (1988)) with fourth order lexicons. On the other
hand Yoshinaka and Kanazawa (2005) have explored the expressivity of lex-
icalized ACGs. They have exhibited a non-semilinear string language with
third order parse structures and an NP-complete string language with fourth
order parse structures. (Salvati (2005) gave an example of an NP-complete
language with third order parse structures and a first order lexicon).

The present work addresses the problem of the expressivity of ACGs in
a particular case. We show that the class of languages defined by second or-
der string ACGs is the same as the class of languages defined as outputs of
deterministic tree walking transducers (DTWT) (Aho and Ullman (1971)).
Together with the results of de Groote and Pogodalla (2004) and Weir (1992),
this result proves that the generative power of second order string ACGs is ex-
actly the same as the generative power of linear context free rewriting sys-
tems. This furthermore shows that second order string ACGs can always be
described with fourth order lexicons. We may nevertheless conjecture that the
use of lexicons of order greater than four may give more compact grammars.

The paper is organized as follows: we first briefly define the linear
λ-
calculus and ACGs in section 11.2. In section 11.3, we use the correspon-
dence between proofs of linear logic and linear λ-terms to relate sub-formulae
of a type α with sub-terms of terms of type α. Section 11.4 introduces h-
reduction, the reduction used by the DTWTs which encode second order
string ACGs. Section 11.5 presents the encoding of second order string ACGs
with DTWTs. Finally we conclude and outline future work in section 11.6.

11.2 Definitions

Given a finite set of atomic types $\mathcal{A}$, we define $\mathcal{T}_\mathcal{A}$, the set of linear applica-
tive types built on $\mathcal{A}$ with the following grammar:

$$\mathcal{T}_\mathcal{A} ::= \mathcal{A} | (\mathcal{T}_\mathcal{A} \to \mathcal{T}_\mathcal{A})$$

If $a_1, \ldots, a_n$ are elements of $\mathcal{T}_\mathcal{A}$ and $\alpha \in \mathcal{A}$ we will write $(a_1, \ldots, a_n) \to \alpha$ the
type $(a_1 \to (\cdots (a_n \to \alpha) \cdots))$. The order of the type $\alpha$, $\text{ord}(\alpha)$, is 1 if $\alpha$
is atomic (i.e. $\alpha \in \mathcal{A}$), and $\text{ord}(\alpha \to \beta) = \max(\text{ord}(\alpha) + 1, \text{ord}(\beta))$.

Higher-order signatures are triples $(\mathcal{C}, \mathcal{A}, \tau)$ where $\mathcal{C}$ is a finite set of con-
stants, $\mathcal{A}$ is a finite set of atomic types and $\tau$ is a function from $\mathcal{C}$ to $\mathcal{T}_\mathcal{A}$. The order of a signature $(\mathcal{C}, \mathcal{A}, \tau)$ is \(\max\{\text{ord}(\tau(a)) | a \in \mathcal{C}\}\). Given a higher-order
signature $\Sigma = (\mathcal{C}, \mathcal{A}, \tau)$ we will denote $\mathcal{A}$ by $\mathcal{A}_\Sigma$, $\mathcal{C}$ by $\mathcal{C}_\Sigma$, $\tau$ by $\tau_\Sigma$ and $\mathcal{T}_\mathcal{A}$ by
$\mathcal{T}_\Sigma$; if $\tau_\Sigma(a) = (a_1, \ldots, a_n) \to \alpha$, then the arity of $a \in \mathcal{C}_\Sigma$ is $n$, it will be noted
$\rho_\Sigma(a)$ or $\rho(a)$ (when $\Sigma$ is clear from the context).

A higher-order signature $\Sigma$ is said to be a string signature if $\mathcal{A}_\Sigma = \{\ast\}$,
$\# \in \mathcal{C}_\Sigma$, $\tau_\Sigma(\#) = \ast$ and for all $a \in \mathcal{C}_\Sigma \setminus \{\#\}$, $\tau_\Sigma(a) = (\ast \to \ast)$.

We are now going to define the set of linear $\lambda$-terms built on a signature
Given a higher-order signature $\Sigma$ and $\alpha \in T_\Sigma$, we assume that we are given an infinite enumerable set of variables $x^\alpha$, $y^\alpha$, $z^\alpha$, $\ldots$, $\Lambda^\alpha_\Sigma$ the set of linear $\lambda$-terms of type $\alpha$ built on $\Sigma$ is the smallest set verifying:

1. if $a \in C_\Sigma$ and $\tau_\Sigma(a) = \alpha$ then $a \in \Lambda^\alpha_\Sigma$
2. $x^\alpha \in \Lambda^\alpha_\Sigma$
3. if $t_1 \in \Lambda^{\alpha \rightarrow \beta}_\Sigma$, $t_2 \in \Lambda^\beta_\Sigma$ and $\text{FV}(t_1) \cap \text{FV}(t_2) = \emptyset$ then $(t_1 t_2) \in \Lambda^\alpha_\Sigma$
4. if $t \in \Lambda^\alpha_\Sigma$, $x^\alpha \in \text{FV}(t)$ then $\lambda x^\alpha.t \in \Lambda^{\alpha \rightarrow \beta}_\Sigma$

The set $\Lambda_\Sigma$ denotes $\bigcup_{\alpha \in T_\Sigma} \Lambda^\alpha_\Sigma$. Linear $\lambda$-terms are linear because variables may occur free at most once in them and that whenever $\lambda x^\alpha.t$ is a linear $\lambda$-term, $x^\alpha$ has exactly one free occurrence in $t$. Moreover, whenever $t \in \Lambda^\alpha_\Sigma \cap \Lambda^\beta_\Sigma$ then $\alpha = \beta$, i.e. every linear $\lambda$-term has a unique type in a given signature $\Sigma$.

We may, when it is not relevant, strip the typing annotation from the variables. We will write $\lambda x_1 \ldots x_n.t$ for the term $\lambda x_1 \ldots \lambda x_n.t$ and $t_0 t_1 \ldots t_n$ for $(t_0(t_1 \ldots t_n))$. Given a list of indices $S = [i_1, \ldots, i_n]$, we will write $\lambda x_{S}^{i_1} \ldots x_{S}^{i_n}.t$ the term $\lambda x_{i_1} \ldots \lambda x_{i_n}.t$, $t_{i_0 i_1} \ldots t_{i_0 i_1 \ldots i_n}$ and $c_{S}^{j} t$ the term $c_{i_1} \ldots c_{i_n}(t \ldots )$ when for all $j \in [1, n]$, $c_{i_j}$ has type $\ast \rightarrow \ast$. In particular, $\lambda x_{S}^{i_1} \ldots x_{S}^{i_n}.t$, $t_{i_0 i_1} \ldots t_{i_0 i_1 \ldots i_n}$ and $c_{S}^{j} t$ may be used when $S = [1, \ldots, n]$.

Given a string signature $\Sigma$, strings will be represented by the closed terms of type $\ast$. For example, the term $c_1(c_2(\ldots c_n(\ldots )))$ represents the string $c_1 \ldots c_n$; given $w$, a string built on $C_\Sigma$, $[w]$ will denote the term of $\Lambda^\ast_\Sigma$ which is in normal form and represents $w$.

To define the sub-terms of $t \in \Lambda_\Sigma$, we follow Huet (1997) and consider them as pairs $(C[], t')$ (where $C[]$ is a context, i.e. a term with a hole) such that $t \equiv C[t']$. The set of sub-terms of $t$ is denoted by $S_t$. In particular, we define $S_t^\beta$ to be $[(C[]), v] \in S_t^\beta \subset \Lambda^\beta_\Sigma$ if $x$ is free in $t$, we note $C_{i_1}[]$ the context such that $C_{i_1}[x] = t$ and $x$ is not free in $C_{i_1}[]$. Remark that since $t$ is linear $C_{i_1}[]$ is always defined.

We say that a term $t$ is in long form if for all $(C[], t') \in S_t^\beta$ either $t' = \lambda x.t''$ or $C[] = C'[][t''/t']$. Every term can be put in long form by $\eta$-expansion, therefore if $t$ is the long form of $t'$, then $t \xrightarrow{\eta} t'$. When a term is in long form, all its possible arguments are abstracted by a $\lambda$-abstraction. For example, the term $x^\ast \rightarrow y^\ast$, which is not in long form, can be applied to an argument of type $\ast$; in long form, this term becomes $\lambda y^\ast.x^\ast \rightarrow y^\ast$. The possibility of applying it to a term of type $\ast$ is syntactically represented by the $\lambda$-abstraction. A term is in long normal form (inf for short) if it is both in $\beta$-normal form and in long form. The set $\text{Inf}_\Sigma^\alpha$ (resp. $\text{Clnf}_\Sigma^\alpha$) represents the set of terms of $\Lambda^\alpha_\Sigma$ in $\text{Inf}$ (resp.

---

1. Given a $\lambda$-term $t$, we will write $\text{FV}(t)$ to denote the set of its free variables.
the closed terms of $\Lambda^a_2$ in $\text{lnf}$). In the sequel of the paper we only deal with terms in long form; thus each time we will write $\lambda x^a_1.t$, $x^a_1$ or $at^a_s$, we will implicitly make the assumption that $t$, $x^a_1$ or $at^a_s$ has an atomic type.

We define homomorphisms between the higher-order signatures $\Sigma_i$ and $\Sigma_2$ to be pairs $(f, g)$ such that $f$ is a mapping from $T_{\Sigma_i}$ to $T_{\Sigma_2}$, and $g$ is a mapping from $\Lambda_{\Sigma_1}$ to $\Lambda_{\Sigma_2}$, and verifying:

1. if $a \in A_{\Sigma_1}$ then $f(a)$ is in $T_{\Sigma_2}$, otherwise, $f(\alpha \rightarrow \beta) = f(\alpha) \rightarrow f(\beta)$
2. for all $a \in C_{\Sigma_1}$ such that $\tau_{\Sigma_1}(a) = \alpha$, $g(a)$ is in $\text{clnf}^{f(a)}$.
3. $g(t) = x^{f(a)}$
4. $g((t_1)_2) = g(t_1)g(t_2)$
5. $g(\lambda x^a t) = \lambda x^{f(a)}(g(t))$

One can easily check that whenever $t \in \Lambda^a_2$, $g(t) \in \Lambda^{f(a)}_2$. In general, given a homomorphism $L = (f, g)$, we will write indistinctly $L(a)$ for $f(a)$ and $L(t)$ for $g(t)$. The order of $L$ is max(ord($L(a)$))$a \in A_{\Sigma_1}$).

An ACG (de Groote (2001)) is a 4-tuple $(\Sigma_1, \Sigma_2, L, S)$ such that:

1. $\Sigma_1$ is a higher-order signature, the abstract vocabulary.
2. $\Sigma_2$ is a higher-order signature, the object vocabulary.
3. $L$ is a homomorphism from $\Sigma_1$ to $\Sigma_2$, the lexicon.
4. $S \in A_{\Sigma_1}$.

An abstract constant (resp. object constant) is an element of $C_{\Sigma_1}$ (resp. $C_{\Sigma_2}$), an abstract type (resp. object type) is an element of $T_{\Sigma_1}$ (resp. $T_{\Sigma_2}$). Given an abstract constant $a$, $L(a)$ is called the realization of $a$.

An ACG $G = (\Sigma_1, \Sigma_2, L, S)$ defines two languages:

1. the abstract language: $\mathcal{A}_1(G) = \text{clnf}^{\Sigma_1}_a$
2. the object language: $\mathcal{O}(G) = \{v \in \text{clnf}^{\Sigma_2}_a | \exists r \in \mathcal{A}_1(G). v =_\eta L(r)\}$

An ACG $G = (\Sigma_1, \Sigma_2, L, S)$ is said to be a string ACG if $\Sigma_2$ is a string signature and $L(S) = \ast$. The order of an ACG is the order of its abstract signature.

11.3 Path in types, active sub-terms and active variables

We assume that we are given a signature $\Sigma$ and that all the types and all the terms used in this section are built on that signature.

A linear $\lambda$-term $t \in \text{lnf}_2$ represents, via the Curry-Howard isomorphism, a cut-free proof of $\alpha$ in the Intuitionistic Implicative and Exponential Linear Logic. This correspondence leads to a natural relation between sub-formulae of $\alpha$ and sub-terms of $t$. This section presents this relation which will play a central role in our encoding.

The sub-formulae of a type will be designated by means of paths. A path $\pi = i_1 \cdot i_2 \cdots i_{n-1} \cdot i_n$ is a possibly empty sequence of strictly positive integers;
Given a path \( \pi \), we define \( p + \pi \) as: \( p + \pi = \begin{cases} \bullet & \text{if } \pi = \bullet \\ (p + k) \cdot \pi' & \text{if } \pi = k \cdot \pi' \end{cases} \)

Given \( t \in \text{Inf}_n \), we define two particular subsets of \( S_t \), the set of active sub-terms, \( \mathcal{AT}_t \), and the set of active variables, \( \mathcal{AV}_t \). The sets \( \mathcal{AT}_t \) and \( \mathcal{AV}_t \) are defined as the smallest sets satisfying:

1. \( ([], t) \in \mathcal{AT}_t \)
2. if \( (C[\lambda a_i \ldots a_n], t') \in \mathcal{AT}_t \), then for all \( i \in [1, n] \),
   \( (C[\lambda a_i \ldots a_n], t_i) \in \mathcal{AV}_t \)
3. if \( (C[\lambda t_1 \ldots t_n], x) \in \mathcal{AV}_t \), then for all \( i \in [1, n] \),
   \( (C[xt_1 \ldots t_{i-1} \ldots t_n], t_i) \in \mathcal{AT}_t \)

If a term \( t \) can be applied to \( n \) arguments, then, given \( t_1, \ldots, t_n \) terms in \( \text{Inf} \), during the \( \beta \)-reduction of \( t_1 \ldots t_n \) the active variables of \( t \) will eventually substituted by a term during \( \beta \)-reduction and the residuals of the active sub-terms of \( t \) will eventually become the argument of a reduct. On the other hand, the variables of \( t \) which are not active will never be substituted and the sub-terms of \( t \) which are not active will never be the argument of a reduct.

We can now define two mutually recursive functions \( \mathcal{AT}_t \) and \( \mathcal{AV}_t \) respectively from \( \mathcal{AT}_t \) onto \( \mathcal{AT}_t \) and from \( \mathcal{AV}_t \) onto \( \mathcal{AV}_t \):

1. \( \mathcal{AT}_t([], t) = \bullet \)
2. if \( \mathcal{AT}_t(C[\lambda a_i \ldots a_n], t') = \pi \) then for all \( i \in [1, n] \),
   \( \mathcal{AV}_t(C[\lambda a_i \ldots a_n], t_i) = \pi \cdot i \)
3. if \( \mathcal{AV}_t(C[\lambda t_1 \ldots t_n], x) = \pi \) then for all \( i \in [1, n] \),
   \( \mathcal{AT}_t(C[xt_1 \ldots t_{i-1} \ldots t_n], t_i) = \pi \cdot i \)

One can easily check that \( \mathcal{AT}_t(C[\lambda], \pi) = \pi \) (resp. \( \mathcal{AV}_t(C[\lambda], x) = \pi \)) implies that the type of \( \nu \) (resp. \( \pi \)) is the type designated (in the obvious way) by \( \pi \) in \( \alpha \).

The functions \( \mathcal{AT}_t \) and \( \mathcal{AV}_t \) are bijections whose converse is \( \mathcal{P}_t \):

1. \( \mathcal{P}_t(\bullet) = ([], t) \)
2. \( \mathcal{P}_t(\pi \cdot i) = \begin{cases} (C[\lambda a_i \ldots a_n], t_i) \text{ if } \mathcal{P}_t(\pi) = (C[\lambda a_i \ldots a_n], t') \\ (C[xt_1 \ldots t_{i-1} \ldots t_n], t_i) \text{ if } \mathcal{P}_t(\pi) = (C[\lambda t_1 \ldots t_n], x) \end{cases} \)
For all \((C[, t']) \in \mathcal{AT}_t\) (resp. \((C[, x]) \in \mathcal{AV}_t\)) it is straightforward that \(P_t(\mathcal{AT}_t(C[, t'])) = (C[, t'])\) (resp. \(P_t(\mathcal{AV}_t(C[, x])) = (C[, x])\)); and that for all \(\pi \in \mathcal{P}_u\) (resp. \(\pi \in \mathcal{P}_u\)), \(\mathcal{AT}_t(P_t(\pi)) = \pi\) (resp. \(\mathcal{AV}_t(P_t(\pi)) = \pi\)).

11.4 \(h\)-reduction

The DTWTs which encode second order string ACGs perform the normalization of the realization of abstract terms. They use a particular reduction strategy, \(h\)-reduction, which is related to head linear reduction (Danos and Regnier (2004)).

This reduction strategy is only defined for a particular class of \(\lambda\)-terms. Firstly, these \(\lambda\)-terms have to be built on a string signature \(\Sigma\); secondly, they have a particular form. To describe this form, we need first define \(N_\Sigma^\alpha \subseteq N_\Sigma^\beta\) (\(N_\Sigma = \bigcup_{\alpha \in \mathcal{N}_\Sigma} N_\Sigma^\alpha\)) as:

\[
N_\Sigma^\alpha := \lnf_\Sigma^\alpha | (N_\Sigma^{\beta = \alpha})
\]

Then, the set of terms we are interested in are the \(HT\)-terms defined by the following grammar:

\[
\mathcal{HT} := N_\Sigma^\alpha | c\mathcal{HT} | (\lambda x_1^{\alpha_1} \ldots x_n^{\alpha_n}. \mathcal{HT}) N_\Sigma^{\beta_1} \ldots N_\Sigma^{\beta_m}
\]

where \(c \in C_\Sigma\). Every \(HT\)-term is in \(N_\Sigma^\alpha\) and is of the form:

\[
(\lambda x_1^{\alpha_1}, c_T^\alpha, \ldots (\lambda x_n^{\alpha_n}, c_T^\beta(x, j^Q)) v_{S_1} \ldots) v_{S_2}^\alpha
\]

so that \(S_i \cap S_j \neq \emptyset\) implies that \(i = j, v_k\) (with \(k \in \bigcup_{\alpha \in [1, m]} S_i\)) and \(t_q\) (with \(q \in Q\)) are elements of \(N_\Sigma\).

Given a \(HT\)-term,

\[
t = (\lambda x_1^{\alpha_1}, c_T^{\alpha_1}, \ldots (\lambda x_n^{\alpha_n}, c_T^{\beta_1}(x, j^Q)) v_{S_1}^\alpha \ldots) v_{S_2}^\alpha
\]

we say that \(t\) \(h\)-contracts to \(t'\) (noted \(t \rightarrow_h t'\)) if

\[
t' = (\lambda x_1^{\alpha_1}, c_T^{\alpha_1}, \ldots (\lambda x_n^{\alpha_n}, c_T^{\beta_1}(v, j^Q)) v_{S_1}^\alpha \ldots) v_{S_2}^\alpha
\]

where \(S_1' = S_1 \setminus \{j\}\). It is a routine to check that \(t =_h t'\), that \(t'\) is also a \(HT\)-term and that the normal form of \(t\) can be obtained in a finite number of \(h\)-contractions. The reflexive and transitive closure of \(\rightarrow_h\), \(h\)-reduction, will be written \(\rightsquigarrow_h\).

Given \(G = (\Sigma_1, \Sigma_2, S, L)\) a second order string ACG, and \(u \in \text{clnf}_\Sigma^\alpha\), we are going to see how \(h\)-contraction normalizes \(L(u)\). The determinism of \(\rightarrow_h\)

allows one to predict statically (i.e. without performing the reduction) which sub-term of \(L(u)\) will be substituted to a given bound variable in \(L(u)\) during \(h\)-reduction. This prediction is based on the notions of replaceable variables and unsafe terms introduced by Böhm and Dezani-Ciancaglini (1975). Replaceable variables and unsafe terms of \(u\) belong to \(S_{L(u)}\) and will be respectively denoted by \(\mathcal{RV}_u\) and \(\mathcal{UT}_u\).
If \((C[, a) \in S_v\) and \((C[, x) \in \mathcal{R}V_{L(a)}, then \((L(C)[C[, x) \in \mathcal{R}V_u; UT_u\) is the smallest set verifying:

1. if \((C[, a) \in S_u\) and \(C[] \neq []\) then \(L(C)[C[, L(a)^w\) \(\in UT_u\)
2. if \((C[, a) \in S_u\) and \((C[, v) \in \mathcal{R}T_{L(a)}\) then \((L(C)[C[, v) \in UT_u\)

The prediction will be given by \(\phi_u\), a bijection between \(\mathcal{R}V_u\) and \(UT_u\).

The definition of \(\phi_u\) relies on few more technical definitions.

Given \((C_a[, a) \in S_u\) such that \(C_a[) = C[[v_1 \ldots v_n,]\), then

\[\var{C[av_1 \ldots v_{i-1}[] \ldots v_n, v_i} \]

is the \(i^{th}\) argument of \((C_a[, a)\). Given \((C_a[, a), (C_b[, b) \in S_u\) we say that \((C_a[, a)\) is the head of the \(i^{th}\) argument of \((C_b[, b)\) if \(C_b[) = C[[v_1 \ldots v_{i-1}(a^w\) \(\ldots v_n,]\) and \(C_a[) = C[[v_1 \ldots v_{i-1}(a^w\) \(\ldots v_n,]\)

Given \((C[, x) \in \mathcal{R}V_u\), we now define \(\phi_u(C[, x)\). As \((C[, x) \in \mathcal{R}V_u\), we have \((C_a[, a) \in S_u\) and \(C[, x)\) such that \((C[, x) \in \mathcal{R}V_{L(a)}\) and \(C[) = L(C_a)(C[, x)\). Let \(\pi = AV_{L(a)}(C[, x), since \(\pi \in P_{L(a)}, \pi\) is of odd length, and \(\pi = i, \pi'\). Then we have three cases:

1. if \(i \leq \rho_a\) and \(\pi' = \bullet\), then \(\phi_u(C[, x) = (L(C'[), L(t)\) where \(C'[[, t}\) is the \(i^{th}\) argument of \((C_a[, a)\)
2. if \(i \leq \rho_a\) and \(\pi' 
eq \bullet\), then \(\phi_u(C[, x) = (L(C'_b)[C'[[, t\) where \(C'_b[, b)\) is the head of the \(i^{th}\) argument of \((C_a[, a)\) and \(C'[[, t) = P_{L(b)}(\rho_b + \pi')\)
3. if \(i > \rho_a\) then \(\phi_u(C[, x) = (L(C'_b)[C'[[, t\) where \(C'_b[, b)\) is the head of the \(i^{th}\) argument of \((C_a[, a)\) and \(C'[[, t) = P_{L(b)}(k \cdot (t - \rho_a) \cdot \pi')\).

Computing \(\phi_u(C[, x)\) only requires to know about the immediate surrounding of \(a\). This is the reason why the normalization of \(L(u)\) can be performed by a DTWT. To prove the correctness of the prediction of \(\phi_u\) we need the notion of strict residual: given \(t\) and \(t'\) such that \(\rightarrow_h t, t', (C[, v) \in S_t\) and \(C[, v) \in S_t\), we say that \((C[, v)\) is the strict residual of \((C[, v)\) whenever \(C[x_1 \ldots y_n, v) \rightarrow_h C'[x_1 \ldots y_n, v)\) with \(FV(v) = \{y_1, \ldots, y_n\}\) and \(x\) is a fresh variable.

Given \(t\) such that \(L(u) \rightarrow_h t\), we say that \(t\) is predicted by \(\phi_u\) if the two following properties hold:

1. for all \((C[, \lambda_{w^m}(\{x\})^w\) \(\in S_t\) and \(i \in [1, n]\), the fact that

\[\var{C[[\lambda_{w^m}(\{x\})^w, x]}\]

is the strict residual of \((C_n[, x)\) \(\in \mathcal{R}V_{L(a)}\) implies that

\[\var{C[[\lambda_{w^m}(\{x\})^w, v]}\]

is the strict residual of \(\phi_u(C_n[, x)\).

2. for all \((C[[v_n^2, x) \in S_t\), \((C[[v_n^2, x)\) is the strict residual of some \((C'[[[v_n^2, x) \in \mathcal{R}V_u\).
We are now going to show that $h$-reduction preserves the predictions of $\phi_a$.
This will be achieved by using the following technical lemma:

**Lemma 25** Given $(C[[v^a]], x) \in \mathcal{RV}_a$ if we have $\phi_a(C[[v^a]], x) = (C'[\lambda], t')$ then $t' = (\lambda \overrightarrow{x_\phi}, \overrightarrow{w})$ and we have 

$$
\phi_a(C'[\lambda], \overrightarrow{w}, C_{w_\phi}[\lambda], v_k) = (C[xv_1 \ldots v_{k-1}] \ldots v_k, v_k)
$$

**Proof.** This proof only consists in unfolding the definitions. Since $(C[[v^a]], x) \in \mathcal{RV}_a$, we must have $(C_a[\lambda], a) \in S_a$ and $C_a[\lambda]$ such that:

1. $C[[v^a]] = L(C_a)[C_a[[v^a]]]
2. (C_a[[v^a]], x) \in \mathcal{RV}_{L(a)}
3. AV_{L(a)}(C_a[[v^a]], x) = i \cdot \pi$ for some $i$ and $\pi$

There are three different cases depending on $i$ and $\pi$.

**Case 1:** $i \leq \rho_a$ and $\pi = \bullet$: this case is very similar to the following one and is thus left to the reader. It is the only case where $p$ may be different from 0.

**Case 2:** $i \leq \rho_a$ and $\pi \neq \bullet$: by definition if $(C_a[\lambda], b)$ is the head of the $i^{th}$ argument of $(C_a[\lambda], a)$, and if $P_{L(b)}(\rho_b + \pi) = (C'', L_b[w])$ then $C'[\lambda] = L(C_a)[C''[\lambda]]$ and $t' = (\lambda \overrightarrow{w}, w)$. Let’s now suppose that $\pi = m \cdot \pi'$, then we have that $AV_{L(a)}(C''[\lambda], v_k) = (\rho_b + \pi') \cdot k = (\rho_b + m) \cdot \pi' \cdot k$. Therefore, as $\rho_b + m > \rho_b$ and as $(C_a[\lambda], b)$ is the head of the $i^{th}$ argument of $(C_a[\lambda], a)$, we have that $\phi_a((L_b[w], C_{w_\phi}[\lambda]), v_k) = (L(C_a)[C''[\lambda]], u_k)$ where

$$(C_a[\lambda], u_k) = P_{L(a)}(i \cdot (\rho_b + m - \rho_b) \cdot \pi' \cdot k) = P_{L(a)}(i \cdot \pi \cdot k)
$$

But we have that $AV_{L(a)}(C_a[[v^a]], x) = i \cdot \pi$ which implies that

$$(C_a[\lambda], u_k) = P_{L(a)}(i \cdot \pi \cdot k) = (C_a[xv_1 \ldots v_{k-1}] \ldots v_k, v_k).
$$

Finally as $C[\lambda] = L(C_a)[C_a[\lambda]]$ we get the result.

**Case 3:** $i > \rho_a$: this case is similar to the previous one.

**Proposition 26** If $L(u) \xrightarrow{h} t$, then $t$ is predicted by $\phi_u$.

**Proof.** This proof is done by induction on the number of $h$-contraction steps of the reduction. The case where this is zero is a simple application of the definitions. Now let’s suppose that $L(u) \xrightarrow{h} t \xrightarrow{h} t'$, then, by induction hypothesis, $t$ is predicted by $\phi_u$; furthermore, $t$ is a HT-term, thus

$$
t = (\lambda \overrightarrow{x_\phi}, C_T(\ldots(\lambda \overrightarrow{x_\phi}, C_T(x, \overrightarrow{Q}) \overrightarrow{s_1} \ldots)) \overrightarrow{s_i})
$$

and

$$
t' = (\lambda \overrightarrow{x_\phi}, C_T(\ldots(\lambda \overrightarrow{x_\phi}, C_T(x, \overrightarrow{Q}) \overrightarrow{s_1} \ldots)) \overrightarrow{s_i}^1)
$$

with $S_i' = S_i[\{j\}]$. 


Within the two conditions required to obtain that \( t' \) is predicted by \( \phi_s \), only the first one requires more than a straightforward application of the induction hypothesis. There is actually only one subterm of \( t' \) which is problematic: \( v_j \). From the induction hypothesis we know that the subterm corresponding to \( x_j \) in \( t \) is the strict residual of \((C[[\rho^j]], x_j) \in R \mathbb{V}_u\) and that the subterm corresponding to \( v_j \) in \( t \) is the strict residual of \( \phi_h(C[[\rho^j]], x_j) \). Finally the previous lemma allows us to conclude that \( v_j \) fulfills the first condition. □

### 11.5 Encoding second order string ACGs with DTWT

We are now going to show how to encode second order string ACGs into DTWT. We do not follow the standard definition of DTWT as given in Aho and Ullman (1971). Indeed, instead of walking on the parse trees of a context free grammar, the transducers we use walk on linear \( \lambda \)-terms built on a second order signature. But, as these sets of \( \lambda \)-terms are isomorphic to regular sets of trees, the string languages output by our transducers are the same as those of usual DTWT. By abuse, we call our transducers DTWT.

A DTWT is defined as a 6-tuple
\[
A = (\Sigma, D, Q, T, \delta, q_0, q_f)
\]
where \( \Sigma \) is a second order signature; \( D \in \mathcal{A}_\Sigma \); \( Q \) is a finite set of states; \( T \) is a finite set of terminals; \( \delta \), the transition function, is a partial function from \( \mathcal{C}_x \times (Q \setminus \{q_f\}) \) to \( (\text{up}, \text{stay}) \cup (\text{down} \times \mathbb{N}^+) \times Q \times T^* \) where \( \mathbb{N}^+ \) denotes the set of strictly positive natural numbers and \( T^* \) denotes the monoid freely generated by \( T \); \( q_0 \in Q \) is the initial state; and \( q_f \in Q \) is the final state. A configuration of \( A \) is given by \((C[, a, q, s) \) where \( C[a] \in \text{clnf}^D_\Sigma \), \( a \in \mathcal{C}_\Sigma \), \( q \in Q \) and \( s \in T^* \); initial configurations are of the form \(([[v]], a, q_0, e) \) (\( e \) being the empty string) where \( v \in \text{clnf}^D_\Sigma \). The automaton \( A \) defines a move relation, \( \triangleright_A \) (\( \triangleright_A \) is the reflexive transitive closure of \( \triangleright_A \)), between configurations: \((C[, a, q, s) \triangleright_A (C[, a, q', s) \) if \( \delta(a, q) = (q', m, w) \) and one of the following holds:

1. \( m = \text{up} \) and \((C[, a) \) is the head of one of the arguments of \((C[, b) \)
2. \( m = \text{stay} \) and \((C[, b) = (C[, a) \)
3. \( m = \text{down} \) and \((C[, b) \) is the head of the \( i \)th argument of \((C[, a) \)

Given \( v \in \text{clnf}^D_\Sigma \), \( \triangleright_A \) generates \( s \) with \( A \) if

\[
([v], a, q_0, e) \triangleright_A (C[, b, q_f, s).
\]

The language of \( A \), \( L_A \), is \{s | \exists v \in \text{clnf}^D_\Sigma \text{ generates } s\}.

Given a second order string ACG \( G = (\Sigma_1, \Sigma_2, \mathcal{L}, S) \) we are going to build an automaton \( A_G = (\Sigma, D, Q, T, \delta, q_0, q_f) \) such that \( O(G) = \{w \mid w \in L_A\} \). Let \( k_G = \max \{\rho(a) \mid a \in \mathcal{C}_\Sigma\} \), we then define \( \Sigma \) as:

1. \( \mathcal{A}_\Sigma = \mathcal{A}_\Sigma \times [1, k_G] \)
2. $C_\Sigma = C_\Sigma \times \{1, k_\Sigma\}$

3. if $\tau_\Sigma(a) = (a_1, \ldots, a_n) \rightarrow a$ then

$\tau_\Sigma((a, k)) = ((a_1, 1), \ldots, (a_n, n)) \rightarrow (a, k)$. 

Remark that if $v \in \text{clnf}^{(\Sigma)}_{\Sigma}$, then for all $(C)[, (a, j)) \in S_v$, $C[] \neq [\tilde{v}]^{\Sigma}$ implies that $(C)[,(a, j))$ is the head of the $j^{th}$ argument of $(C)[,(b, l)) \in S_v$. Furthermore, given $v = (a, k)(\tilde{v}) \in \text{clnf}^{(\Sigma)}_{\Sigma}$ we note $\tilde{v}$ the term of $\text{clnf}^{(\Sigma)}_{\Sigma}$, such that $\tilde{v} = a^v$. 

Then $D = (S, 1), Q = ([0, k_\Sigma] \times P) \cup \{q_i\}$ where $P = \bigcup_{v \in C_\Sigma}\mathcal{P}_{\mathcal{L}(\alpha)}, q_0 = (0, \bullet); \text{building } \delta \text{ requires some more definitions.}$

Given $(a, k)$ and $(i, \pi)$, the selection path of $(a, k)$ and $(i, \pi)$ is:

$$\pi' = \begin{cases} i \cdot \pi & \text{if } i > 0 \\ \rho_\pi + \pi & \text{if } i = 0 \end{cases}$$

If the selection path of $(a, k)$ and $(i, \pi)$ is in $\mathcal{P}_{\mathcal{L}(\alpha, \Sigma)}$, then we say that $(a, k)$ and $(i, \pi)$ are coherent; $\delta$ will be only defined on coherent pairs of $(a, k)$ and $(i, \pi)$. A configuration $K = (C)[, (a, k), (i, \pi), w]$ is said to be coherent if $(a, k)$ and $(i, \pi)$ are coherent.

If $(a, k)$ and $(i, \pi)$ are coherent and if $\pi'$ is their selection path, then we define the focused term of $(a, k)$ and $(i, \pi)$ as $\mathcal{P}_{\mathcal{L}(\alpha)}(\pi')$. Furthermore, if $(C)[, t)$ is the focused term of $(a, k)$ and $(i, \pi)$ and if $t = \lambda \tilde{v}_{\pi'}^{\Sigma}, c_n(x\tilde{v}_{\pi'}^{\Sigma}), \mathcal{X}_{\Sigma}(\lambda \tilde{v}_{\pi'}^{\Sigma}, c_n(x\tilde{v}_{\pi'}^{\Sigma}), x)$ is called the focused variable of $(a, k)$ and $(i, \pi)$.

If $(a, k)$ and $(i, \pi)$ are coherent then $\delta((a, k), (i, \pi)) = (q, move, w)$ depends on the focused term of $(a, k)$ and $(i, \pi)$, (noted $(C)[, t)$):

1. if $t = \tilde{c}_n^{\Sigma}$ then $q = q_1$, move = stay and $w = c_1 \ldots c_n$
2. if $t = \lambda \tilde{v}_{\pi'}^{\Sigma}, c_n(x\tilde{v}_{\pi'}^{\Sigma}), \mathcal{X}_{\Sigma}(\lambda \tilde{v}_{\pi'}^{\Sigma}, c_n(x\tilde{v}_{\pi'}^{\Sigma}), x) = l \cdot \pi''$ and $l > \rho_\pi$ then $q = (k, (l - \rho_\pi) \cdot \pi''), \text{ move = up and } w = c_1 \ldots c_n$
3. if $t = \lambda \tilde{v}_{\pi'}^{\Sigma}, c_n(x\tilde{v}_{\pi'}^{\Sigma}), \mathcal{X}_{\Sigma}(\lambda \tilde{v}_{\pi'}^{\Sigma}, c_n(x\tilde{v}_{\pi'}^{\Sigma}), x) = l \cdot \pi''$ and $l \leq \rho_\pi$ then $q = (0, \pi''), \text{ move = (down, l) and } w = c_1 \ldots c_n$

We now relate the walk of $A_{\Sigma}$ on $v \in \text{clnf}^{(\Sigma, \Sigma)}_{\Sigma}$ with the $h$-reduction of $\mathcal{L}(\tilde{v})$. To establish this relation we need to show that the transducer computes $\phi_{\Sigma}$. Given a coherent configuration $K = (C)[, (a, k), (i, \pi), w)$, the activated term of $K$ is $(\mathcal{L}(C)'[]), (\mathcal{L}(\tilde{v}_{\pi'}^{\Sigma})$ if $(i, \pi) = (0, \bullet)$ and $\tilde{C}[] = C'[\tilde{v}]^{\Sigma}$, otherwise it is $(\mathcal{L}(C)'[[]], t)$ if $C'[[]], t)$ is the focused term of $(a, k)$ and $(i, \pi); \text{ the activated variable of } K = (\mathcal{L}(C)[[]], x)$ if the focused variable of $(a, k)$ and $(i, \pi)$ is $(C[], x)$. We will show that given $K_1$ and $K_2$ such that $K_1 \downarrow_{A_{\Sigma}} K_2$, if $(C[, x)$ is the activated variable of $K_1$ then $\phi_{\Sigma}(C[], x)$ is the activated term of $K_2$. This property shows that $A_{\Sigma}$ performs the $h$-reduction of $\mathcal{L}(\tilde{v})$ and that if $\mathcal{L}(\tilde{v})$ normalizes to $\tilde{w}$ then, walking on $v$, $A_{\Sigma}$ ends in the final state and outputs $w$. 


**Lemma 27** Given \( \nu = (a, 1)^{\nu_{\rho_0}} \in \text{clfn}_{\Sigma}^{(S, 1)} \) and two coherent configurations \( K_1 \) and \( K_2 \) such that \( \((\nu_{\rho_0}, (a, 1), (0, \bullet), \epsilon) \vdash_{\Lambda_{\Sigma}} K_1 \vdash_{\Lambda_{\Sigma}} K_2 \) if \( (C[], x) \) is the activated variable of \( K_1 \) then \( \phi(\Sigma, x) \) is the activated term of \( K_2 \).

**Proof.** As for the proof of lemma 25, this proof is mainly based on the unfolding of the definitions. We simply compute \( \phi(\Sigma, x) \) and the activated term of \( K_2 \) and then show that they are the same.

We assume that \( K_r = (C_r[], (a_r, k_r), (i_r, \pi_r), w_r) \) with \( r \in [1, 2] \), that \( \pi'_r \) is the selection path of \( K_r \). If \( \text{P}_{\Sigma_{(\pi, r)}}(\pi'_r) = (C_r[], A_\Sigma^{k_r, c_n}(x_v^{\pi'}_{r})) \), then let \( \pi''_r = \text{AV}_{\Sigma_{(\pi, r)}}(C_r[], A_\Sigma^{k_r, c_n}(x_v^{\pi'}_{r}, x)) \) as \( \pi''_r \in \text{P}_{\Sigma_{(\pi, r)}}(\pi'_r) \), we know that \( \pi''_r = i \cdot \pi'' \). We then have three cases:

**Case 1:** if \( i \leq \rho_{a_1} \) and \( \pi'' = \bullet \), then \( \phi(\Sigma, x) = (\Sigma(C[]), \Sigma(t)) \) if \( (C[]), t \) is the \( i \)th argument of \( (C[]), a_1 \). But in that case, we have that \( \delta((a_1, k_1), (i, \pi)), (\text{down}, t, c_1 \ldots c_n) \); thus \( (a_2, k_2) \) is the head of the \( i \)th argument of \( (a_1, k_1) \) and as \( (i_2, \pi_2) = (0, \bullet) \), we obtain, by definition, that the activated term of \( K_2 \) is indeed \( (\Sigma(C[]), \Sigma(t)) \).

**Case 2:** if \( i \leq \rho_{a_1} \) and \( \pi'' \neq \bullet \), then \( \phi(\Sigma, x) = (\Sigma(C[]), \Sigma(t)) \) if \( (C[]), \Sigma(t) = \text{P}_{\Sigma_{(\pi, r)}}(\rho_{a_1} + \pi'') \). In that case, we have that \( \delta((a_1, k_1), (i, \pi)), (\text{down}, t, c_1 \ldots c_n) \); thus \( (a_2, k_2) \) is the head of the \( i \)th argument of \( (a_1, k_1) \) which implies that \( (C_2[]), a_2 = (C_1[]), \Sigma(t) \); finally by definition we have that the activated term of \( K_2 \) is \( (\Sigma(C[]), \Sigma(t)) \).

**Case 3:** if \( i > \rho_{a_1} \), then \( \phi(\Sigma, x) = (\Sigma(C[]), \Sigma(t)) \) if \( (C[]), t \) is the head of the \( k \)th argument of \( (C[]), b \) and \( (C[]), t = \text{P}_{\Sigma_{(\pi, r)}}((i - \rho_{a_1}), \pi'') \).

In that case, we have that \( \delta((a_1, k_1), (i, \pi)), (\text{down}, t, c_1 \ldots c_n) \).

The definition leads to the fact that the activated term of \( K_2 \) is \( (\Sigma(C[]), \Sigma(t)) \).

\( \square \)

**Proposition 28** Given \( u \in \text{clfn}_{\Sigma}^{(S, 1)} \), there is a unique \( \nu = (a, 1)^{\nu_{\rho_0}} \in \text{clfn}_{\Sigma}^{(S, 1)} \) such that \( v = u \), and \( (\nu_{\rho_0}, (a, 1), (0, \bullet), \epsilon) \vdash_{\Lambda_{\Sigma}} (C[]), b, q_{f}, w) \) iff \( \Sigma(u) = \Sigma_{\rho_0} [w] \).

**Proof.** The existence and the uniqueness of \( v \) are obvious from the definition of \( \Sigma \). To prove the proposition it suffices to study the walk of \( \Lambda_{\Sigma} \) on \( v \) and the \( h \)-reduction of \( \Sigma(u) \) in parallel: assume that \( K_1 = (\nu_{\rho_0}, (a, 1), (0, \bullet), \epsilon) \), \( t_1 = \Sigma(u) \), \( K_1 \vdash_{\Lambda_{\Sigma}} K_2 \) and \( t_1 \rightarrow_{h} t \) (where \( \rightarrow_{h} \) corresponds to \( k \) steps of \( \Lambda_{\Sigma} \) and \( \rightarrow_{k} \) to \( k \) steps of \( h \)-reduction). The use of the previous lemma and an induction on \( k \) prove that \( t_k \) is of the form

\[ t_k = (\Lambda^{t_{k+1}}(\Sigma_{t_{k+1}}) \ldots (\Lambda^{t_{k+1}}(\Sigma_{t_{k+1}}(x), f_q)) \ldots (\Lambda^{t_{k+1}}(\Sigma_{t_{k+1}}(x)))) \]

if and only if \( K_k = (C[]), (a^2, k_0), (i_k, \pi_k), w_k) \) so that \( w_k = c_{i_k} \ldots c_{i_k} \), if
$(C'_k[\lambda], \lambda \tau_k^T(x_j, \alpha_j)) \in S_n$ (with the obvious $C'_k[]$) is the strict residual of $(C''_k[\lambda], \lambda \tau_k^T(x_j, \alpha_j)) \in S_n$, then $(C''_k[\lambda], \lambda \tau_k^T(x_j, \alpha_j))$ is the activated term of $K_k$ and $(C''_k[\lambda \tau_k^T(\alpha_j), x_j])$ is the activated variable of $K_k$. This allows us to conclude that the walk ends in the configuration $(C[]), b, q_f, w$ iff $L(u) \equiv_{\beta\eta} /w/$. This finally shows that $O(G)$ is indeed equal to $\{/w/ | w \in L_{AG}\}$.

11.6 Conclusions and future work

In this paper, we have proved that the languages defined by second order string ACGs were the same as the output languages of DTWT. From the results of Weir (1992) and de Groote and Pogodalla (2004), we obtain as a corollary that the languages defined by second order string ACGs are exactly the languages defined by LCFRS. Furthermore as, according to de Groote and Pogodalla (2004), LCFRS can be encoded by second order string ACGs with a fourth order lexicons, we obtain that every second order string ACG can be encoded by another one whose lexicon has at most fourth order.

In our next work, we would like to exhibit a direct translation of a second order string ACG into another one with a fourth order lexicon. This would help understanding how relevant the order of the lexicon is. We conjecture that using lexicons of order greater than four may lead to more compact grammars. The problem is to know how compact those grammars can be and if the compaction is important whether it can be used do design large grammars for natural languages.

As the tools we used are general, we think it is possible to prove that any second order ACG can be represented as a second order ACG whose lexicon is at most fourth order. Indeed, the notion of paths and the relations they establish with active sub-terms and active variables do not depend on the problem. The only definition which is dependent of the fact we deal with strings is the definition of $h$-reduction. We nevertheless think that, provided we define a generalized notion of DTWT which would output linear $\lambda$-terms instead of strings, we can show that second order ACGs can be encoded with these generalized DTWTs. It would remain to encode those DTWTs with second order ACGs with a fourth order lexicon to generalize our result. But this last part does not seem too difficult.

The first part seems also feasible since it should be possible to generalize $h$-reduction. Indeed, instead of having a unique variable on which we could make the substitution, the fact that the constants in the term introduce some branching may lead to have several such variables. This would correspond on the generalized DTWTs to the fact that when it would output a branching constant the transducer should duplicate its head in order to have one head to generate each argument of that constant.
Finally this work may lead to the definition of an abstract machine for second order ACGs. Such a machine would be valuable to study the problem of parsing second order ACGs and give insights on the strategies that can be implemented for those grammars. Furthermore, as such a machine would have a language made of linear \(\lambda\)-terms, it would be a first step towards the definition of an abstract machine whose language is a set of \(\lambda\)-terms. In Montague style semantics, the problem of generation mainly consists in parsing languages of \(\lambda\)-terms. We would then obtain a valuable tool to study the problem of generation in that setting.

References


Sidewards without copying

EDWARD P. STABLER

Abstract

A traditional movement step relates a single source position to a single c-commanding target position, and never moves an argument to another argument position. But head movement involves non-c-command relations, and control relates two argument positions that are not always in a c-command relation. Special mechanisms could be invoked for these things, but a different strategy slightly generalizes movement and enforces certain fundamental symmetries observed by all movements to block over-generation. This paper defines a class of ‘sideward movement grammars’ (smmg) with such symmetries, with example applications to adjunct control and head movement. These grammars allow copying, but the question of whether to copy is completely independent of the question of whether to allow sideward movement. Furthermore, since these grammars distinguish complement attachments from others, a simple CED-like constraint can block extractions from specifiers and adjuncts except in the exceptional circumstance of adjunct control. smmg definable languages are all pmcfg definable, and hence are efficiently recognizable.

Keywords  Parsing, grammar, syntax

12.1 Introduction

One of the most basic properties of human language is its simple, recursive, layered character in which similar structure is iterated, sometimes with special variations at the top, matrix level and at the deepest levels:

```
<table>
<thead>
<tr>
<th>Does Alice know that</th>
<th>Bob thinks that</th>
<th>Carol says</th>
<th>you like her?</th>
</tr>
</thead>
<tbody>
<tr>
<td>3</td>
<td>2</td>
<td>1</td>
<td>0</td>
</tr>
</tbody>
</table>
```

Certain kinds of recursive symmetry in languages allow the ‘pumping lemmas’ which have been valuable diagnostics of the availability of certain kinds of grammars. A regular grammar for a language is only possible when the
language has a simple symmetry of this kind; context free grammars have a weaker requirement, and so on through the hierarchy of multiple context free languages (Seki et al., 1991), etc.

Many descriptions of human languages involve rearranging constituents. In grammars with movements, how is the structure of each ‘layer’ affected? This fundamental question is a topic of active study. In early transformational grammars, a set of base structures is generated and then transformed into surface structures, as in the following example (with e and t unpronounced):

\[ [I [know [e [I [e [saw [who]]]]]]] \rightarrow [I [know [who [I [t [saw [t]]]]]]]. \]

The sequences of positions related by movement in these accounts are not random. Among other things, landing sites of movement do not disrupt layer structure too much (‘structure preservation’, ‘shape conservation’), and when an element moves through several clauses, it never moves from a high position in a lower clause to a lower position in a higher clause (cf. the ‘ban on improper movement’ ‘chain uniformity’, ‘level embedding’). So in effect, the hierarchy of each layer of phrase structure is respected in sequences of movements too, another reflection of the basic invariants mentioned at the outset.

Some recent grammars compose generation and transformation steps,\(^1\) so transformations are, in effect, executed as soon as requisite structure is built, reducing the need for revising completed structure:

1. \([\text{saw}] + [\text{who}] \xrightarrow{\text{merge}} [\text{saw} \ [\text{who}]]\]
2. \([\text{saw} \ [\text{who}]] + [I] \xrightarrow{\text{merge}} [I \ [\text{saw} \ [\text{who}]]]\]
3. \([I \ [\text{saw} \ [\text{who}]]] \xrightarrow{\text{move}} [\text{who} \ [I \ [\text{saw} \ [\text{who}]]]]\]
4. \([\text{know}] + [\text{who} \ [I \ [\text{saw} \ [\text{who}]]]] \xrightarrow{\text{merge}} [\text{know} \ [\text{who} \ [I \ [\text{saw} \ [\text{who}]]]]]\]
5. \([\text{know} \ [\text{who} \ [I \ [\text{saw} \ [\text{who}]]]]] + [I] \xrightarrow{\text{merge}} [I \ [\text{know} \ [\text{who} \ [I \ [\text{saw} \ [\text{who}]]]]]]\]

But step 3 shows who being copied and deleted, revising the structure built by step 2. One response is to say that the syntax simply copies the earlier structure (perhaps only adding a link, a pointer to the embedded who), and then a post-syntax “spellout” process determines which copies to pronounce. This pushes the changes to completed structure out of the syntax, by invoking a “spellout” process that is sensitive to much of the same structure that syntactic operations are sensitive to. When two processes seem to be sensitive to the same structure it is a natural hunch that they are really the same process. Adopting this perspective instead, we could then say that the depiction of the derivation 1-5 is slightly misleading: when who is introduced in step 1, it satisfies a requirement of the verb but is not actually placed in complement position. Rather, it is held out to be placed at the left edge of the embedded

clause. This strategy for (not postponing but) eliminating a kind of structural
revision is formalized in mgs (Stabler and Keenan, 2003, Frey and Gärtner,
do not ban improper movements.

Now consider the co-indexed elements in sentences like these:

\[ \text{He} \text{, tries } [e_i \text{ to succeed}] \]
\[ \text{He} \text{, laughs } [\text{before } e_i \text{ eating}] \]

These ‘obligatory control’ (OC) relations have enough in common with
movement to suggest a uniform treatment (Hornstein, 2006, 2001, 1999,
movement so that a subject can move to another subject position even out
of an adjunct as in the latter example, the rest of the phrasal construction
can remain completely standard. But such movements between unconnected
structures must be restricted to avoid unwanted movements, like these for
example:

*John, likes \( t_i \)
*The cook they, like tried \( t_i \) to make them
*John, persuaded Mary \( t_i \) to make them
*John,’s friends prefer \( t_i \) to behave himself

One critique of movement analyses of control wonders, if sideways move-
ment is allowed, what rules out sideward movement from complements gener-
ally (Landau, 2003, p.477). In the present account, the status and restrictions
on sideward movement will be clear: sideward movement from complements
is impossible.

Another kind of problem is posed by head movements like this:

\[- \text{an} + \text{[ustedes [habl- [español]]]} \rightarrow \text{[habl-an] [ustedes [habl- [español]]]} \]

If we say \( x \) c-commands \( y \) in a tree iff a sister of \( x \) dominates \( y \), then \text{habl} does
not c-command its original position. Adapting a proposal from Nunes (2001)
and Hornstein (2001), in analogy to phrasal movement, we can compute this
result without surgery by keeping the head \text{habl}- out of its projection so that is
available for attachment to the appropriate affix. But the indicated assembly
of the head and affix with the rest of the projection is more complicated than any
of the other (merge,move) rules, looking suspiciously \textit{ad hoc}. An alternative
is to, in effect, allow the head to move before it projects its structure. This
yields essentially the same result, but by allowing the head to simply move to
another projection, allows the construction of the phrase and the selection of
that phrase to be completely standard. But obviously this step needs to bring
some analog of the traditional head movement constraint (HMC):

*be -s he have be -en making tortillas
Conventional movements relate source constituents with targets that command them. In mgs, the same effect is achieved by keeping the sources separate from the target while they wait for their final licensed positions. In this setting, the needed generalization simply allows new, ‘disconnected’ elements to be inserted into an expression. With this generalization of expressions, we need only one feature-checking operation, *merge*. We define ‘sideward movement grammars’ (SMMGs) in this way. To avoid overgeneralization, we impose a specifier island constraint (SpIC) and also impose a generalized ban on improper movements. Since all phrases other than the matrix clause are either complements or specifiers, SpIC allows extracted phrases to enter a derivation only through complements, though as explained below this constraint is weaker than usual because a complement can be remnant-moved to a specifier without freezing any of its moving elements.

Formal antecedents include tree adjoining grammar (Joshi and Schabes, 1997) and especially the variants proposed for scrambling (Rambow et al., 2001, Rambow, 1994, Kallmeyer, 1999), certain elaborations of pregroup grammars (Stabler, 2004a, Casadio and Lambek, 2002, Buszkowski, 2001), and the minimalist grammars (mgs) already mentioned. The derivations in these formalisms all extend and simplify complexes of possibly discontinuous constituents. But none of them enforces the ban on improper movements, and none of them defines the same class of languages as SMMGs. SMMG languages are not all MCRG definable, but they are all PMCGR-definable (Seki et al., 1991) and hence are polynomially parsable. We conjecture that all PMCGR languages are SMMG definable too.

### 12.2 Sideward movement grammars

Let $\Sigma$ be a finite vocabulary, associated with phonetic and semantic properties. The empty sequence is $\epsilon$. Head movement will be triggered by a morphological property that we indicate with hyphens: a preceding hyphen $-$ indicates that a lexical head is a suffix; a following hyphen $-$ indicates a prefix; and the affix $s$ can be empty.

A set of syntactic features $\mathbb{F}$ is partitioned into 2 basic kinds: properties $\mathbb{F}$ and requirements $\mathbb{F}$. Properties $\mathbb{F}$ are either persistent $-f$ or not $-$. Requirements $\mathbb{F}$: some simply require agreement $+f$, others trigger overt movement $+f$, and others trigger overt movement and also leave a copy $+f$. As in MGS, we use the types $\mathbb{T} = \{\text{;}\text{;}\}$ to indicate lexical and derived expressions, respectively. The projections $\mathbb{P} = \Sigma^* \times \mathbb{T} \times \mathbb{F}$. The expressions $\mathbb{E} = \mathbb{P} \times \wp(\mathbb{P})$. Consider, e.g., the expression

\[(\text{loves}:-v, \{\text{Mary}:-\text{focus}, \text{who}:-\text{case}$-$wh\}).\]

To reduce clutter, we often omit some braces and parentheses,
loves:-v, Mary:-focus, who:-case -wh.

With this simpler notation, remember that the head of an expression comes first, and the order of remaining elements (if any) is irrelevant.

A lexicon is a finite subset of $\Sigma^* \times \{\cdot\} \times (\{+\mathbb{F}^+\} \times \{\cdot\}) \times \{\emptyset\}$ with a designated ‘start’ category $f$. A lexical item has category $f$ iff its first property is $-f$ or $-\cdot$. $f$ comp-selects $g$ iff there as a lexical item with category $f$ whose first requirement is $+g$ or $+\cdot$ or $+g$. A cycle is a sequence $f_0 \ldots f_n$ such that $f_0$ is the start category, $f_{i-1}$ comp-selects $f_i$ (all $0 < i \leq n$), and no feature appears twice. $f$ cycle-selects $g$ iff $f$ precedes $g$ in a cycle. A lexicon is proper iff whenever $-f$ precedes $-g$ in any lexical item, some lexical item containing $-f$ has category $c$ and some lexical item containing $-g$ has category $d$, where $d$ cycle-selects $c$. With this constraint on lexicons, (Proper), we can remain neutral about whether human languages have a universal, fixed clausal structure. A grammar is given by a proper lexicon, generating the structures in the closure of lexicon with respect to the fixed structure building rules. A completed structure is one containing only one syntactic feature, the start category $f$. The string language is the set of yields of those completed structures.

There are two structure building relations, ins and merge. The partial binary function $\text{ins}$ applies to pairs of expressions $((p, S), (q, T))$ only if (i) either $(q, T)$ is lexical or $S = \emptyset$, and (ii) $\text{match}(p, q)$ is defined. Its value is given by $\text{ins}((p, S), (q, T)) = (p, S \cup \{q\} \cup T)$. Condition (i) is our version of SpIC, mentioned above.

The relation $\text{merge} \subset E \times E$ applies to $(p, S)$ only if there is a unique $q \in S$ such that $\text{match}(p, q)$ is defined. Then it takes as value $\text{merge}(p, S \cup \{q\}) = (r, (S - q) \cup T)$ for each $\text{match}(p, q) = (r, T)$. The uniqueness condition on application of this function is our version of the shortest move constraint (SMC).

The relation $\text{match} \subset P \times P \times E$ is given as follows, where $s, t \in \Sigma^*$ are not marked with an initial or final hyphen to trigger head movement, $\alpha, \beta, \gamma \in \mathbb{F}^+$, $\delta \in \mathbb{F}^+$, and $\cdot \in T$. 

We present some examples to illustrate these mechanisms and set the stage for introducing sideward movement.

**Example 1: Basics.** In the derivation tree on the left, the leaves are lexical items; The binary branches represent applications of insert, and the unary branches, applications of merge.

\[
\begin{array}{c|c|c}
\text{Overt movement:} & p & q \\
\hline
s::fa & t^{-f} & s:ta,\emptyset & \text{saturated complement} \\
 s::fa & t^{-f} & ts:ta,\emptyset & \text{saturated specifier} \\
 s::fa & t^{-t} & s:ta,(t:\delta) & \text{moving, unsaturated projection} \\
 s::fa & t^{-f} & s:ta,\emptyset & \text{final use of -f} \\
 s::fa & t^{-f} & ts:ta,\emptyset & \text{final use of -f} \\
 s::fa & t^{-t} & s:ta,(t:\delta) & \text{moving, unsaturated projection} \\
 s::fa & t^{-f} & s:ta,\emptyset & \text{moving with -f} \\
 s::fa & t^{-f} & ts:ta,\emptyset & \text{final use of -f} \\
 s::fa & t^{-t} & s:ta,(t:-f) & \text{moving, unsaturated projection} \\
 s::fa & t^{-f} & s:ta,\emptyset & \text{moving with -f} \\
 s::fa & t^{-t} & s:ta,(t:-f) & \text{check non-persistent -t} \\
 s::fa & t^{-t} & s:ta,(t:-f) & \text{final use of -f} \\
 s::fa & t^{-f} & s:ta,(t:-f) & \text{moving with -f} \\
 s::fa & t^{-t} & s:ta,(t:-f) & \text{moving with -f} \\
\end{array}
\]

\[
\begin{array}{c|c|c}
\text{Covert movement:} & p & q \\
\hline
s::fa & t^{-t} & s:ta,\emptyset & \text{saturated complement (i)} \\
 s::fa & t^{-t} & ts:ta,\emptyset & \text{saturated specifier (ii)} \\
 s::fa & t^{-t} & s:ta,(t:\delta) & \text{moving, unsaturated projection (iii)} \\
 s::fa & t^{-t} & s:ta,\emptyset & \text{moving, unsaturated projection (iv)} \\
 s::fa & t^{-t} & s:ta,\emptyset & \text{moving, unsaturated projection (v)} \\
 s::fa & t^{-t} & s:ta,\emptyset & \text{moving with -f (vi)} \\
 s::fa & t^{-t} & s:ta,\emptyset & \text{moving with -f (vii)} \\
 s::fa & t^{-t} & s:ta,\emptyset & \text{moving, unsaturated projection (viii)} \\
 s::fa & t^{-t} & s:ta,\emptyset & \text{final use of -f (ix)} \\
 s::fa & t^{-t} & s:ta,\emptyset & \text{moving with -f (x)} \\
\end{array}
\]

\[
\begin{array}{c|c|c}
\text{Copy movement:} & p & q \\
\hline
s::fa & t^{-t} & s:ta,\emptyset & \text{saturated complement (xi)} \\
 s::fa & t^{-t} & ts:ta,\emptyset & \text{saturated specifier (xii)} \\
 s::fa & t^{-t} & s:ta,(t:\delta) & \text{moving (xiii)} \\
 s::fa & t^{-t} & s:ta,(t:\delta) & \text{moving (xiv)} \\
 s::fa & t^{-t} & s:ta,\emptyset & \text{moving (xv)} \\
 s::fa & t^{-t} & s:ta,\emptyset & \text{moving to complement (xvi)} \\
 s::fa & t^{-t} & s:ta,\emptyset & \text{moving to specifier (xvii)} \\
 s::fa & t^{-t} & s:ta,\emptyset & \text{moving to complement (xviii)} \\
\end{array}
\]
Note that since insert applies to introduce a projection that can be merged, and the derivation greedily checks features at the earliest possible moment, there is a merge immediately above each insert step. The additional unary branches represent ‘external merge’ steps: these are the steps that are traditionally called ‘movements’. The tree on the right shows the corresponding conventional X-bar structure. It is not difficult to translate the derivations shown here into more traditional depictions like this.  

**Example 2: Obligatory control into a complement.** One idea about obligatory control is that there is a special unpronounced pronoun PRO which, unlike other pronouns, either does not need case or else needs some special kind of case that infinitival tense can assign. But Hornstein argues that the PRO positions can be the empty positions left by movement, as in:

```
he tries to succeed:-C
\[ \epsilon::+T^-C, he \text{ tries to succeed}:-T \]
```

This derivation is checking the categorial D feature of [he] twice (and then checking its case feature in a higher clausal position, in conformity with Proper). Hornstein suggests that really it is θ-features getting checked twice in constructions like this. (And there have been suggestions that categorial...

---

2This translation can be done automatically. See the implementations at http://www.linguistics.ucla.edu/people/stabler/epssw.htm.
features generally should be replaced by appropriate complexes of more basic features: \( \theta \)-features etc.) For present purposes, the simple analysis above provides a suitable starting point.

**Example 3: Obligatory control into an adjunct.** There are many interesting questions about adjunction, but for present purposes it suffices to adopt a treatment that allows it to be category-preserving, iterable, optional, and opaque to extraction. These properties can be obtained by introducing an empty category to host the adjunct; for clausal adjuncts of noun phrases we use \( \varepsilon : + N \leftrightarrow N \wedge N - N \), and for prepositional modifiers of \( v \) we can use: 
\[ \varepsilon : + v \leftrightarrow P + v - v \], as in:

\[
\text{he laughs before he eats: -C}
\]

\[
\varepsilon : + T \leftrightarrow C \text{, he laughs before he eats: -T}
\]

The fact that \([\text{before he eats}]\) is a specifier is indicated by the non-lexical status of the selector \( \varepsilon : + P + v - v, \text{laughs:} - v, \text{he:} - k, \text{before he eats:} - P \). Since SpIC blocks any extraction from specifiers, we do not need to separately stipulate that adjuncts are islands. So if we introduce right and left \( X \)-adjuncts of \( Y \) with lexical items of the form \( \varepsilon : + X \leftrightarrow Y + X - X \), or \( \varepsilon : + X + Y - X \), respectively (or with any processes that yields similar structure), we get the desired properties.
for adjuncts: optionality, iterability, and opacity to extraction. This sets the stage for the special treatment of adjunct control.

Since the proposed treatment of adjuncts makes them opaque to extraction, while the proposed treatment of control makes it an extraction relation, we should not get control into adjuncts, but we do:

he, laughs before eating

Hornstein notices that a slight tweak on our mechanisms can let this kind of case through without allowing other kinds of adjunct extractions. Roughly, if we derive the modifier \([\text{before } e_i \text{ eating, } \{\text{he}\}]\) which wants to attach to a \(v\), and then we derive a \(v\) that is looking for a \(D\), we can allow \([\text{he}]\) to ‘move sideways’ onto the \(v\) before inserting it into the derivation. This step can be presented in logicians’ style, as the inference from the expressions above the line to the one below:

\[
\begin{align*}
\text{before eating} &: -P, [\text{he} : -D-k] \\
\epsilon &: +v+P+\text{-}v, \emptyset \\
\text{laughs before eating} &: -v, [\text{he} : -D-k]
\end{align*}
\]

We express this step more generally as follows. In a grammar that contains left \(X\)-adjuncts of \(Y\), that is, it has some

\[r = e::+X+Y+X-X\]

we extend the (ins) relation so that it also applies to \(((p, [a]), (q, S))\) in the exceptional case where \(p\) and \(q\) can be chained together by \(r\), using \(a\) as follows:

\[
\begin{align*}
\text{match}(q, a) &= (b, T), \\
\text{match}(r, b) &= (c, U), \\
\text{match}(c, p) &= (e, V), \text{ and} \\
\text{match}(e, f) &= (g, W) \text{ for } f \in U.
\end{align*}
\]

Notice that the adjoining element \(r\) is introduced in the second step to have its 3 initial features checked in sequence. In this special case, let

\[\text{ins}((p, S), (q, T)) = (g, S \cup T \cup (U - \{f\}) \cup V \cup W).\]

Control into right \(X\)-adjuncts of \(Y\) can be defined similarly, using the lexical item \(\ell = e::+X+Y-X\), checking its 2 initial features in sequence. With this extension, we obtain:
Example 4: Head movement is similar to adjunct control in relating constituents that do not c-command each other, but, unlike control, we want just the phonetic parts of the heads to move while their projections are developed in their original positions. Nevertheless, there is an application of the side-ward movement idea that avoids splitting all phrases kept into triples so that the head can be separate when the phrase is complete, as was done in Stabler (2001).

We extend match so that, when the category of $-s:\alpha$ is comp-selected by $t:\beta$ and $t-s$ is morphologically well-formed,

$$p \quad q \quad \text{match}(p, q)$$

\begin{align*}
-s:\alpha & \quad t:\beta & \quad e:\alpha, [t-s:\beta] & \quad \text{suffixed left adjoins lower head} \\
s:\alpha & \quad t:\beta & \quad e:\alpha, [s-t:\beta] & \quad \text{prefix right adjoins lower head}
\end{align*}

And then, when match($q, p$) is defined by one of (i-xviii) we bring the adjunction up:

$$p \quad q \quad \text{match}(p, q)$$

\begin{align*}
p & \quad q & \quad q.(p) & \quad \text{higher head promoted}
\end{align*}

With these extensions, we get derivations like the following:
12.3 Expressive power and recognition complexity

Previous studies have shown that head movement, though it may seem like a small thing in informal presentations, allows the definition of non-context free patterns even when there is no phrasal movement in the grammar. But the translation from mgs to mcgs defined by Michaelis (2001) is easily adapted to show that smmg grammars without copying all define mcg definable languages. There are various theory-internal arguments for copying in grammar, and various ways to implement them (Stabler, 2004b). See for example Nunes (2001) and Kobele (2006) for some empirical arguments in support of rather powerful copy operations. The addition of copy features makes it easy to define non-semilinear languages like $a^{2^n}$, but a straightforward extension of Michaelis’s translation to these cases shows that they are pmcdefinable, and hence polynomially recognizable.

12.4 Conclusions

This paper does not attempt to resolve the controversy over whether movement analyses of obligatory control are empirically well-motivated (Landau, 2003, Boeckx and Hornstein, 2004), but provides a formalization of some parts of these ideas that can be rigorously studied.
Although smmg\$s can be regarded as extending mgs, notice that they differ in a number of significant respects: (1) smmg\$s extend the domain of movement just slightly to offer tightly constrained treatments of obligatory control and head movement. Future work may find ways to make these constraints more general and natural. And there are regularities in the definition of match that should allow a more elegant statement. (2) mgs are bound by SMC, while smmg\$s also are required to respect SpIC and Proper, and future work may provide further additions. (3) To handle head movement, mgs require either extra rules for head movement (Michaelis, 2001) or else one of the approaches mentioned in the introduction. smmg\$s allow head movement with a simple mechanism analogous to the sideward mechanisms used for control. (4) mgs have no copy operation, and while none of the analyses above depend on it, smmg\$s allow copying. That is, we have presented a treatment of sideward movement that does not rely in any way on the copy theory of movement for its appeal. In the present setting, sideward movement is a natural option not because we already have operations on copies, but because we already have operations on moving phrases (the original phonetic materials, not copies). smmg\$s are naturally extended to allow copying though, setting the stage for studying proposals about overt copying (Boeckx et al., 2005, for example) – unfortunately beyond the scope of this short report. All the mechanisms proposed here are obtained in the well-understood and feasible space of pmcfg-definable languages.

References


English prepositional passives in HPSG

JESSE TSENG

Abstract
This paper provides a detailed syntactic description of English prepositional passives (also known as “pseudopassives”) and discusses their formal treatment in HPSG. The empirical overview includes a discussion of the familiar (but unformalizable) notion of semantic cohesiveness, as well as new observations about the possibility of elements intervening between V and P. Two formal approaches to the syntactic aspects of the problem are then outlined and compared—one relying on lexical rules, the other taking advantage of HPSG’s capacity to express constraints on constructions.

Keywords Pseudopassives, prepositions, adjuncts, HPSG, lexical rules, constructions

13.1 Empirical observations
English has an exceptionally rich variety of preposition stranding phenomena, perhaps the most striking of which is the prepositional passive—the possibility of passivizing the object of a preposition instead of the direct object of a verb.

(24) a. You can rely [on David] to do get the job done.
    b. David, can be relied on ti to get the job done.

Here the NP David, initially the complement of on, is realized as the subject of the passive verb relied, leaving the preposition behind.\(^1\)

It is often suggested that the underlined verb and preposition in this construction form a kind of “compound”, an intuitive notion that is open to many

\(^1\)I will occasionally use the symbol “\(t\)” to mark the “deep” position of the passive subject, in cases where there might be ambiguity. This is deliberately reminiscent of NP-trace in transformational analyses, but here it should be understood only as a expository device with no theoretical strings attached.
formal interpretations. I will begin by presenting some attempts to characterize the phenomenon in semantic terms, before turning to the syntactic aspects of the structure, which will be the main focus of the rest of the paper.

13.1.1 Semantic cohesion

One semantic approach that dates back at least to the classic descriptions of Poutsma and Jespersen is the idea that the prepositional passive is possible if there is a high degree of “cohesion” between V and P. Variants of this position can be found in modern grammars (e.g., Quirk et al., 1985) and in theoretical work on preposition stranding phenomena (see Hornstein and Weinberg (1981), who propose that V and P must form a “natural predicate” or a “possible semantic word”). The most accessible indicator of semantic cohesion is the possibility of replacing the V+P sequence by a single-word synonym:

(25) David can be relied on to get the job done.

But this criterion can easily be shown to be unreliable indicator of passivizability. In particular, many perfectly natural-sounding prepositional passives have no appropriate corresponding one-word synonym (26). Conversely, replacing an ordinary passivized transitive verb by a synonymous V+P combination often produces a degraded result (27), although as discussed below, I do not consider such prepositional passives to be syntactically ill-formed.

(26) That bridge is too low to be undersailed/underpassed/undergone.

(27) a. I was approached (by a complete stranger).
   b. ??I was moved/come/walked towards (by a stranger).

It has also been observed that V+P sequences with abstract, transferred, or metaphorical meaning are more cohesive (i.e., they are more likely to allow the prepositional passive) than concrete, literal uses of the same sequence:

(28) a. An acceptable compromise was finally arrived at.
   b. ??A picturesque mountain village was finally arrived at.

The difference in acceptability between these two examples must be due to non-syntactic factors, since under normal assumptions they receive identical syntactic analyses. Similarly, semantically non-compositional and idiomatic V+P combinations should be expected to be more cohesive and give rise to good prepositional passives, and this is generally the case. The usefulness of these observations for the current study is limited, however, because prepositional passives formed from fully compositional, concrete V+P sequences are generally grammatical, too. They may have relatively degraded acceptability as isolated examples, like (28b), or they can be completely unproblematic, like (26).
Other authors have attempted to approach the prepositional passive by looking at the semantic properties of the targeted oblique NP. Bolinger (1977, 1978) proposes that this NP can become the passive subject if it refers to a strongly “affected” patient. As Riddle and Sheintuch (1983) note, no satisfactory definition is provided for this “dangerously wide” notion, and it is easy to find examples of grammatical prepositional passives where affectedness is not involved. Their own functional account (relying on the notion of “role prominence”) is equally vague.²

Cohesion and affectedness are of course gradient properties, and they can no doubt be decomposed into more primitive, interacting factors. For example, modality, tense, and negation have all been found to influence the acceptability of the prepositional passive. Furthermore, examples that are dubious in isolation can always be improved with an enlarged context.

In this paper I make the simplifying assumption that any V+PP combination can give rise to a syntactically well-formed prepositional passive. The acceptability of the resulting structure, however, is conditioned by non-syntactic restrictions that are not well enough understood to be incorporated into a formal analysis. Existing semantic accounts may be intuitively appealing but they lack a precise, empirical basis, especially if we take into account the predominant role of context. It is also clear that more or less idiosyncratic lexical properties associated with specific V+P combinations are a major determining factor in the acceptability of the prepositional passive; I will abstract away from such considerations in the following, primarily syntactic discussion.

13.1.2 Adjacency
A directly observable sign that V and P form a kind of “compound” in prepositional passive constructions is the fact that the insertion of adverbs and other material between V and P is generally disallowed, whereas various kinds of intervening elements are possible between V and PP in the corresponding active structure:

(29) We rely increasingly [on David] ~ *David is relied increasingly on.

This evidence suggests a constraint on syntactic structure and/or surface word order. ³ This restriction could be formalized by introducing a word order constraint requiring V and P to be adjacent in the passive case, but for various reasons this approach would be inadequate.

²They themselves note that it is “impossible to offer an algorithm for determining what causes some entity or concept to be viewed as role prominent.”

³Note that preposition stranding by extraction is much freer in this respect (although there are restrictions, probably of a prosodic nature):

(i) We rely increasingly [on David] ~ Who do we rely increasingly on?
The specifier *right*, for instance, is possible with some spatial and temporal Ps:

(30) Mr. Cellophane may be looked right through, walked right by and never acknowledged by those who have the audacity to suppose that they cannot be looked right through.

Similar examples can be found with other PP specifiers (*straight, clear, etc*.), so this is not a lexical idiosyncrasy of the word *right*. And in fact, in cases like these, where the preposition has clear relational meaning, a wider variety of modifiers can (quite marginally) appear between V and P in the passive:

(31) The bridge must be ***walked halfway across, sailed completely under, or driven quickly over*** (for the point to be awarded).

Unlike PP specifiers, which must appear immediately the left of P, the placement of the modifiers in this example is clearly “non-optimal”, since they could also appear after P instead, leaving V and P adjacent. There are obviously semantic factors at work here that need to be explored further. From a syntactic point of view, adjacency of V and P is not a strict requirement; I will assume in this paper, in particular, that P can combine with a specifier or a modifier to its left.

Nominal elements can also separate V and P in the prepositional passive. It is well known that passives can be formed from some fixed expressions and light verb constructions containing a bare N or full NP:

(32) a. We were opened fire on, made fools of, paid attention to, taken unfair advantage of.

b. *That product can’t be made a profit from.*

The commonly accepted assumption is that ordinary NP objects cannot appear between V and P, and the prepositional passive is indeed quite bad in most examples of this type:

(33) Samuel explained a complicated theorem to David. *David was explained a complicated theorem to.*

A richer context can significantly improve such examples, however, and some examples of the same structure [V NP P] are unexpectedly good even with minimal context:

---

4This example is from a letter to the editor of the *Bradford Telegraph & Argus* (5 June 2003), referring to lyrics from a song: “Mr. Cellophane shoulda been my name, ’cause you can look right though me, walk right by me, and never know I’m there.”

5Again, the contrast with extraction constructions is striking:

(i) Samuel explained a complicated theorem to David. *Who did Samuel explain a complicated theorem to?*
(34) ?[To be whispered such dirty innuendoes about] would be enough to
drive anyone crazy.

According to Bolinger (1977, 1978), the underlined direct object in this sen-
tence functions as part of the predicate, and the passive subject (left unex-
pressed here) is strongly “affected” by being whispered-dirty-innuendoes-
about. Another proposal by Ziv and Sheintuch (1981) requires such inter-
vening direct objects to be “non-referential”. This is a reasonable characteri-
zation of the idiomatic examples in (32), but in order to accommodate cases
like (34), the authors are forced to broaden the commonly understood notion
of non-referentiality considerably, and to admit that it is “not a discrete prop-
erty”. In the end, the acceptability of this kind of prepositional passive (and
of all prepositional passives, for that matter) depends primarily on context,
and on usage and frequency effects associated with specific lexical items (or
combinations of lexical items).

What is clear is that there can be no strict structural constr aint against the
presence of a direct object in the prepositional passive construction (e.g., an
adjacency condition). We can also demonstrate that the ungrammaticality of
the prepositional passive in cases like (33) is not due to the linear position
of the direct object (between V and P). Even if the object is realized in a
different position, making V and P adjacent, the prepositional passive does
not become more acceptable. On the contrary, the passive examples below,
with V adjacent to P, are worse than example (33) above, with intervening
NP:

(35) a. Samuel explained to David [a fantastically complicated theorem
about the price of cheese]. (heavy NP shift)
   b. *David, was explained to ti [a fantastically complicated theorem
about the price of cheese].

(36) a. the theorem that Samuel explained to David / Which theorem did
Samuel explain to David? (extraction)
   b. *the theorem that David, was explained to ti / *Which theorem was
David, explained to ti?

Furthermore, in cases like (32), where an intervening direct object is un-
problematic, there appears to be a sort of “anti-adjacency” condition on V and
P. Although the direct object can be realized in various positions in the active
voice, in the prepositional passive it must appear between V and P:

(37) a. the unfair advantage that [they took of us] / How much advantage
did they take of us? (extraction)
   b. *the unfair advantage that [we were taken of] / *How much advan-
tage were we taken of?
(38) a. We could make from this product [the kinds of profits that no one
    has ever dreamed of] (heavy NP shift)
    b. *This product, could be made from it [the kinds of profits that no
        one has ever dreamed of].

Based on these observations, I make the following assumptions for the re-
mainder of this paper:

- The prepositional passive is syntactically compatible with the presence of
  a direct object.
- The direct object must be realized in its canonical position between V and
  P.
- The acceptability of the prepositional passive is ultimately determined by
  non-syntactic factors that (for now) resist formalization.

To my knowledge, only one other kind of element can intervene between
V and P in the prepositional passive: When a phrasal verb is involved, its
particle must appear in this position:

(39) a. This situation will simply have to be put up with it.
    b. The loss in speed can be made up for by an increase in volume.

This is unsurprising, given the strong restrictions on particle placement in
English. In the active voice, the particle must be realized closest to the verb
(in the absence of a direct object); this constraint continues to apply in the
passive.6

13.1.3 Further observations

Most of the examples given so far involve passive subjects originating in com-
plement PPs, but it is clear that prepositional passives can also be formed from
[V + adjunct PP] structures:

(40) a. This bed has not been slept in.
    b. David always takes that seat in the corner because he hates being
       sat next to.

The most common sources are temporal and locative modifiers, but we also
find other PPs, like instrumental with-phrases. Again, I will not attempt to
identify or formalize the relevant semantic and lexical constraints. For the

---

6Examples of verbs selecting a particle, a direct object, and a PP at the same time show
that the relative order of the particle and the object remains the same in the active and in the
prepositional passive:

(i) a. They kept an eye out for David. ~ David was kept an eye out for.
    b. *They kept out an eye for David. ~ *David was kept out an eye for.

---
moment, I simply note that the possibility of passivizing out of adjuncts constitutes a crucial difference between the prepositional passive and the ordinary passive.  

We might also wonder if there is any difference between the two passives in terms of their morphological effects, given that they target different (but overlapping) sets of verbs. In particular, the prepositional passive applies to intransitive verbs like sleep or go, and to prepositional verbs like rely, which never undergo ordinary passivization. For verbs that do participate in both types of passivization, we might ask if two distinct morphological operations can be identified. In fact, there is no evidence for this. In every case, the same participial form is used in both constructions:

(41)  
(a) The pilot flew the airplane under the bridge.  \(\rightsquigglyarrow\) The airplane was flown under the bridge. (ordinary passive)  
(b) The pilot flew under the bridge.  \(\rightsquigglyarrow\) The bridge was flown under t. (prepositional passive)

It is not the case that (say) a strong participle flown is used for the ordinary passive, while a weak form *flied is used in the prepositional passive. Both passives require a form of the verb identical to the past participle.\(^8\)

Finally, I briefly discuss the formation of deverbal adjectives from passive V+P sequences:

(42)  
(a) our effective, relied-upon marketing strategy  
(b) a first novel from an as yet unheard-of author

This is sometimes taken as an additional argument for "cohesion" between V and P in the prepositional passive. For example, Hornstein and Weinberg (1981) use it to motivate the semantic notion of "possible word". It is unclear, however, what these adjectives can tell us about the passive structures they derive from, since they are evidently subject to additional constraints. Not all prepositional passives can be used to derive prenominal adjectives:

(43)  
(a) ??a sailed-under bridge, *a sat-beside grouch

---

\(^7\)NP adjuncts, for any number of reasons, cannot passivize like direct objects:  
(i) The children slept three hours.  \(\rightsquigglyarrow\) *Three hours were slept (by the children).

\(^8\)One apparent counterexample is the following pair:  
(i)  
(a) They laid the sleeping child on the rug.  \(\rightsquigglyarrow\) The child was laid on the rug.  
(b) The child lay on the rug.  \(\rightsquigglyarrow\) *The rug was lain/laid on by the child.

Here it looks as if a single verb can have a special participial form lain in the prepositional passive. But in fact two distinct verbs are involved in these examples: transitive lay (with past participle laid) vs intransitive lie (past participle lain/laid). This pair causes confusion and hesitation for most speakers in the past and perfect. It is safe to say, however, that no speaker merges the two into a single verb while maintaining distinct passive forms as in (41).
Some of these examples could be improved with more context, but they all clearly have a degraded status with respect to their fully acceptable verbal counterparts. This is particularly true for the examples with an NP or particle intervening between V and P. The data suggest strongly that adjectival derivation is not a truly productive process, but is more or less lexicalized on a case by case basis. This could perhaps be accounted for with a usage-based model, but I will not pursue this idea any further here.

13.2 Implications for an HPSG analysis

13.2.1 Modularity

The normal passive construction (with the direct object NP “promoted” to subject) is standardly handled as a lexical phenomenon in HPSG, either using a lexical rule deriving the passive participle from an active base verb (Pollard and Sag, 1987), or by assuming an underspecified verbal lexeme that can be resolved to either an active or a passive form with the appropriate linking constraints (Davis and Koenig, 2000).

A number of other approaches can be imagined and technically implemented within the framework, although they have never been seriously explored. For example, passive verbs could have the same syntactic valence as active verbs, if new syntactic combination schemas were added that realized their \textit{comps} element (direct object) in subject position and their \textit{subj} element as a coindexed \textit{by}-phrase. This analysis assumes a different division of labor between lexical information and syntactic operations, but it does not seem to present any advantages in return for the additional complexity it introduces.

A more radical solution would be to approximate the old transformational analysis within HPSG. A recent trend in the framework (most fully developed in Ginzburg and Sag (2001)) is the use of constructional constraints, a departure from the original emphasis (perhaps over-emphasis) on lexical descriptions as the driving force behind syntactic derivation. One characteristic of the constructional approach is a reliance on nonbranching (“head-only”) syntactic rules. Such rules can potentially be used to encode arbitrarily abstract syntactic operations, from a simple change of bar level (e.g., \textit{X}\textsubscript{0} to \textit{XP}), to a coercion of one syntactic category into another (e.g., \textit{S} to \textit{NP}), or in our case, even the transformation of an active clause into a passive clause.

This last proposal would be soundly rejected by linguists working in HPSG, for violating various well-motivated locality and modularity principles. In particular, a syntactic rule should not be able to refer to or arbitrarily modify the phonological, morphological, or internal syntactic structure of the
The proposed non-branching passive transformation rule would have to do all of the above. The problem is that these locality and modularity principles cannot be formally enforced in HPSG; they have the status of conceptual guidelines that responsible practitioners of the theory agree to follow by convention. Of course, this is a fundamental issue that is relevant for all grammatical frameworks, and rarely addressed. But the “all-in-one” sign-based architecture that constitutes the principal strength of HPSG, also makes it particularly easy to fall afoul of these basic principles. In the case of the passive, a constraint requiring non-branching rules to leave the phonology and morphology values unchanged would be enough to invalidate the undesirable transformational analysis. But this is nothing more than an artificial stipulation, covering only a small subset of cases, and the more general theoretical question remains.

13.2.2 Adjunct analyses

For the ordinary passive construction, a strictly lexical analysis is available, because it only needs to refer to the subject and direct object, both of which are present in the lexically defined “argument structure” (encoded in the arg-st list). The fact that PP adjuncts can be involved in prepositional passives (recall the examples in (40)), however, makes a lexical approach to the phenomenon more problematic. This is because information about the identity of eventual adjuncts is not normally available at the lexical level, at least not according to the original assumptions of HPSG. A technical work-around to this problem is possible, in the form of the dependents list of Bouma et al. (2001). This list, whose value is defined as the lexical arg-st extended by zero or more (underspecified) adjuncts, was introduced in order to allow a uniformly head-driven analysis of extraction from complement and adjunct positions.

This result is made possible basically by treating some adjuncts as complements, from a syntactic point of view. This reverses the direction of selection in adjunct structures: The head now selects these adjuncts, in complete contrast to the treatment of adjuncts in Pollard and Sag (1994). This move potentially introduces significant problems for semantic composition. Levine (2003) discusses a problem involving adjuncts scoping over coordinated structures, and argues for a return to the earlier HPSG approach, with adjuncts introduced at the appropriate places in the syntactic derivation (perhaps as empty elements, if they are extracted). Sag (2005) offers a response, requiring modifications to the proposal by Bouma et al. but maintaining the treatment of certain adjuncts as elements selected lexically by the head (and a traceless analysis of extraction).
13.2.3 Prepositional passive: lexical approach

In light of this active controversy, any phenomenon involving adjuncts can be approached in two very different ways in HPSG. At first sight, the adjuncts-as-complements approach seems more appropriate for the prepositional passive, precisely because it targets complement and adjunct PPs in the same way. The lexical rule in Figure 1 takes as input a base form (active voice) verb with a PP on its OBJECT list and outputs a passive participle with a passive specification custom-built to generate the prepositional passive: The first element on OBJECT is the subject, followed optionally by a particle or a direct object. The direct object, if present, is constrained to be canonical, to account for the data in (37–38) above. (An extracted or extraposed/shifted phrases would correspond instead to a non-canonical subtype of synsem.) The crucial operation in this lexical rule is the replacement of a saturated PP (complement or adjunct) in the input by a unsaturated P in the output description. The unrealized complement of the preposition is coindexed with the passive subject NP, and the original subject is optionally realized in a by-phrase, as in the ordinary passive construction.

The complexity and ad hoc nature of this rule is perhaps forgivable, given the highly exceptional status of the phenomenon it models. On the other hand, the proposal fails to capture what is common to the prepositional passive and the ordinary passive. In fact, most aspects of the prepositional passive could be handled by the existing rule for the ordinary passive, which already provides a mechanism for: promoting a non-subject NP to subject position, demoting the subject NP to an optional by-phrase, and ensuring the appropriate morphological effects (identical for both kinds of passive, as confirmed in §13.1.3). For this to work, the NP complement of P must be made available directly on the OBJECT list of the base verb (by applying argument raising, familiar from HPSG analyses of French and German non-finite constructions) so

---

9 This simplified formulation does not accommodate structures containing both a particle and an object (recall fn. 6).

10 E.g., Hinrichs and Nakazawa (1994) and Abeillé et al. (1998).
it can be input to the general passive rule. But this means introducing a systematic ambiguity between the sublists ⟨PP⟩ and ⟨P, NP⟩ in the deps value of the active form of the verb, giving rise to two structures:

\[(44) \quad \begin{array}{ll}
   \text{a.} & \text{VP} \\
   & \text{[rely on David] [deps } \langle \text{NP, PP} \rangle ] \\
   \text{b.} & \text{VP} \\
   & \text{[rely on NP, P, NP] [deps } \langle \text{NP, P, NP} \rangle ]
\end{array} \]

The unwanted analysis (44b) should be blocked, although we need this version of the verb rely in order to generate the prepositional passive was relied on. One straightforward way to achieve this would be to add the specification non-canonical to the second NP element on the verb’s deps list. This would make it impossible for it to be realized as a complement, as in (44b), but we would still have spurious ambiguity in extraction constructions (where the NP is in fact non-canonical). A more adequate solution would be to enrich the hierarchy of synsem subtypes to encode the syntactic function of the corresponding phrase. This would then allow us to state the appropriate constraint (e.g., “¬comps-synsem”).¹¹

This analysis of the prepositional passive is still incomplete, because the insertion of intervening modifiers between V and P must be restricted; recall the discussion of example (29). The lexical operations proposed so far manipulate the deps list, a rather abstract level of representation that cannot be used to express constraints on surface word order. The required constraints therefore have to be formulated separately.

13.2.4 Prepositional passive: syntactic approach

A more radical treatment can be developed for the prepositional passive by combining the earlier HPSG approach to adjuncts (as unselected elements introduced in the syntax) and the more recent trend of constructional analysis.

Figure 2 sketches a special head-adjunct rule that can be used to construct the adjunct-based examples in (40). As in an ordinary head-adjunct phrase, semantic composition is handled via mon selection. But this rule is extraordinary in that it requires the adjunct to be comp-s-unsaturated, and it specifies the coindexation of the unrealized complement of P and the as-yet-unrealized subject of the resulting VP. The rule also imposes special constraints on the head daughter. The sign type core-vp is defined to be compatible with a bare V, or a combination of V with a particle and/or a direct object. In other words, as soon as a verb combines with a non-nominal complement or any kind of

¹¹This can be thought of as a very weak kind of inside-out constraint (as used in LFG, and reinterpreted for HPSG by Koenig (1999)).
modifier, the resulting phrase is no longer a core-vp. This constraint (which constitutes a minor violation of locality principles) determines what can and cannot intervene between V and P in the prepositional passive, as discussed in §13.1.2. The negative constraint on the head daughter’s comps list and the empty slash specification ensure that the particle and object (if any) are actually realized within the core-vp.\textsuperscript{12} There is no particular constraint on the internal structure of the adjunct daughter: It can be either a bare preposition, or a phrasal projection including a specifier or a modifier.

A number of additional details need to be worked out; in particular, some aspects of passivization (e.g., morphological effects) must still be dealt with at the lexical level. It should also be noted that a similar special version of the head-complement rule is needed for prepositional passives involving PP complements, although it is possible to factor out the shared aspects of the two constructional rules; this is precisely the advantage of the hierarchical approach to constructions in HPSG. These preliminary observations suggest that the constructional treatment provides a more satisfactory account of the phenomenon than the lexical approach. Additional questions for further work include a comparison with the prepositional passive in Scandinavian, and a search for similar phenomena anywhere outside of the Germanic family.

References


\textsuperscript{12}This presupposes a return to syntactic slash amalgamation, as in the original HPSG Non-local Feature Principle.


Linearization of affine abstract categorial grammars

Ryo Yoshinaka

Abstract

The abstract categorial grammar (ACG) is a grammar formalism based on linear lambda calculus. It is natural to ask how the expressive power of ACGs increases when we relax the linearity constraint on the formalism. This paper introduces the notion of affine ACGs by extending the definition of original ACGs, and presents a procedure for converting a given affine ACG into a linear ACG whose language is exactly the set of linear $\lambda$-terms generated by the original affine ACG.

Keywords: Abstract categorial grammars, generative capacity, lambda calculus, context-free tree grammars, linear context-free rewriting systems, multiple context-free grammars

14.1 Introduction

De Groote (2001) has introduced abstract categorial grammars (ACGs), in which both lexical entries of the grammar as well as grammatical combinations of them are represented by simply typed linear $\lambda$-terms. While the linearity constraint on grammatical combinations is thought to be reasonable, admitting non-linear $\lambda$-terms as lexical entries may allow ACGs to describe linguistic phenomena in a more natural and concise fashion.

On the other hand, de Groote and Pogodalla (2003, 2004) have shown that a variety of context-free formalisms, namely, context-free grammars, linear¹ context-free tree grammars (linear CFTGs)² and linear context-free rewriting systems, are encodeable by ACGs.

¹This paper lets the term “linearity” mean non-duplication and non-deletion. Thus “linear CFTGs” means non-duplicating non-deleting CFTGs here, though usually “linear CFTGs” means non-duplicating CFTGs.

²See also Kanazawa and Yoshinaka (2005) for complete proof of encodability of linear...
ing systems (LCFRSs), is encoded by ACGs in straightforward ways. In this sense, ACGs can be thought of as a generalization of those grammar formalisms. The linearity constraint in those formalisms matches that of the ACG formalism.

Concerning those grammar formalisms, it is known that the expressive power does not change when the linearity constraint is relaxed to just non-duplication, allowing deleting operations. Seki et al. (1991) have shown the equivalence between LCFRSs and multiple context-free grammars (MCFGs), which correspond to the relaxed version of LCFRSs that may have deleting operations. Fujiyoshi (2005) has established the equivalence between linear monadic CFTGs and non-duplicating monadic CFTGs. Fisher’s result (Fisher, 1968a,b) is rather general. He has shown that the string IO-languages generated by general CFTGs coincide with the string IO-languages generated by non-deleting CFTGs.

Along this line, extending the definition of usual linear ACGs, this paper introduces affine ACGs, which have BCK-λ-terms as their lexical entries, and compares the generative power of linear ACGs and affine ACGs. We present a procedure for converting a given affine ACG into a linear ACG whose language is exactly the set of the linear λ-terms generated by the original ACG. Therefore, affine ACGs are not essentially more expressive than linear ACGs, since strings and trees are usually represented with linear λ-terms.

As linear ACGs encode linear CFTGs and LCFRSs, affine ACGs encode non-duplicating CFTGs and MCFGs in straightforward ways. For such affine ACGs, our linearization method constructs linear ACGs which have the form corresponding to linear CFTGs or LCFRSs. Thus, our result is a generalization of the results we have mentioned above with the exception of Fisher’s, which covers CFTGs involving duplication.

14.2 Preliminaries

14.2.1 Lambda-Terms

Let $\mathcal{A}$ be a finite non-empty set of atomic types. The set $\mathcal{T}(\mathcal{A})$ of types built on $\mathcal{A}$ is defined as the smallest superset of $\mathcal{A}$ such that

- if $\alpha, \beta \in \mathcal{T}(\mathcal{A})$, then $(\alpha \rightarrow \beta) \in \mathcal{T}(\mathcal{A})$.

The order of a type is given by the function $\text{ord} : \mathcal{T}(\mathcal{A}) \rightarrow \mathbb{N}$, where

- $\text{ord}(p) = 1$ for all $p \in \mathcal{A}$,
- $\text{ord}(\alpha \rightarrow \beta) = \max\{\text{ord}(\alpha) + 1, \text{ord}(\beta)\}$.

A higher-order signature $\Sigma$ is a triple $(\mathcal{A}, \mathcal{C}, \tau)$ where $\mathcal{A}$ is a finite non-empty set of atomic types, $\mathcal{C}$ is a finite set of constants, and $\tau$ is a func-
tion from $\mathcal{C}$ to $\mathcal{T}(\mathcal{A})$. The order of the higher-order signature is defined as $\text{ord}(\Sigma) = \max\{\text{ord}(\tau(a)) | a \in \mathcal{C}\}$.

Let $\mathcal{X}$ be a countably infinite set of variables. The set $\Lambda(\Sigma)$ of $\lambda$-terms (terms for short) built upon $\Sigma$ and the type $\hat{T}(M)$ of a term $M \in \Lambda(\Sigma)$ are defined inductively as follows:

- For every $a \in \mathcal{C}$, $a \in \Lambda(\Sigma)$ and $\hat{T}(a) = \tau(a)$.
- For every $x \in \mathcal{X}$ and $\alpha \in \mathcal{T}(\mathcal{A})$, $x^\alpha \in \Lambda(\Sigma)$ and $\hat{T}(x^\alpha) = \alpha$.
- For $M, N \in \Lambda(\Sigma)$, if $\hat{T}(M) = (\alpha \rightarrow \beta)$, $\hat{T}(N) = \gamma$, then $(MN) \in \Lambda(\Sigma)$ and $\hat{T}((MN)) = \beta$.
- For $x \in \mathcal{X}$, $\alpha \in \mathcal{T}(\mathcal{A})$ and $M \in \Lambda(\Sigma)$, $(\lambda x^\alpha.M) \in \Lambda(\Sigma)$ and $\hat{T}((\lambda x^\alpha.M)) = (\alpha \rightarrow \hat{T}(M))$.

For convenience, we simply write $\tau$ instead of $\hat{T}$ and often omit the superscript on a variable if its type is clear from the context. The notions of free variables, closed terms, $\beta$-normal form, $\beta\eta$-normal form, are defined as usual (see Hindley (1997) for instance). A term $M$ is a combinator if $M$ is closed and $M$ contains no constants. A term $M$ is said to be affine if any variable occurs free at most once in every sub-term of $M$. An affine term is said to be linear if every $\lambda$-abstraction binds exactly one occurrence of a variable. The sets of affine and linear terms are respectively denoted by $\Lambda^{\text{aff}}(\Sigma)$ and $\Lambda^{\text{lin}}(\Sigma)$. As usual, let $\rightarrow_{\beta}, =_{\beta}, =_{\beta\eta}$ denote $\beta$-reduction, $\beta$-equality, $\beta\eta$-equality, and $\alpha$-equivalence respectively. $|M|_{\beta}$ and $|M|_{\beta\eta}$ respectively represent the $\beta$-normal form and $\beta\eta$-normal form. We use upper case italic letters $M, N, P, \ldots$ for terms, late lower case italic letters $x, y, z, \ldots$ for variables, middle lower case italic letters $o, p, \ldots$ for atomic types, Greek letters $\alpha, \beta, \ldots$ for types, sanserif $a, A, \ldots$ for constants. We write $a \rightarrow \beta \rightarrow y \rightarrow \delta$ for $(a \rightarrow (\beta \rightarrow (y \rightarrow \delta)))$, $a^\beta \rightarrow \delta$ for $a \rightarrow a \rightarrow a \rightarrow \delta$, $\text{MNQ}$ for $((((M)N)P)Q)$, $\text{Mxyz}$ for $(\lambda x.(\lambda y.(\lambda z.M)))$, and so on.

### 14.2.2 Abstract Categorial Grammars

**Definition 12** For two sets of atomic types $\mathcal{A}_0$ and $\mathcal{A}_1$, a type substitution $\sigma$ is a mapping from $\mathcal{A}_0$ to $\mathcal{T}(\mathcal{A}_1)$, which can be extended homomorphically as

$$\sigma(\alpha \rightarrow \beta) = \sigma(\alpha) \rightarrow \sigma(\beta).$$

For two higher-order signatures $\Sigma_0$ and $\Sigma_1$, a term substitution $\theta$ is a mapping from $\mathcal{C}_0$ to $\Lambda(\Sigma_1)$ such that $\theta(a)$ is closed for all $a \in \mathcal{C}_0$. For two higher-order signatures $\Sigma_0$ and $\Sigma_1$, we say that a type substitution $\sigma : \mathcal{A}_0 \rightarrow \mathcal{T}(\mathcal{A}_1)$ and a term substitution $\theta : \mathcal{C}_0 \rightarrow \Lambda(\Sigma_1)$ are compatible if $\sigma(\tau_0(a)) = \tau_1(\theta(a))$ holds for all $a \in \mathcal{C}_0$. A lexicon from $\Sigma_0$ to $\Sigma_1$ is a compatible pair of a type substitution and a term substitution. A lexicon $\mathcal{L} = \langle \sigma, \theta \rangle$ is affine (linear) if $\theta(a)$ is affine (linear) for all $a \in \mathcal{C}_0$. For a lexicon $\mathcal{L} = \langle \sigma, \theta \rangle$, we define $\hat{\theta}$ as the homomorphic extension of $\theta$ such that $\hat{\theta}(x^\alpha) = x^{\sigma(\alpha)}$. Indeed, $\hat{\theta}(M)$ is...
always a well-typed $\lambda$-term if so is $M$; if $M$ has type $\sigma$, then $\hat{\theta}(M)$ has type $\sigma'(\alpha)$.

Hereafter we identify a lexicon $L = \langle \sigma, \theta \rangle$ with the functions $\sigma$ and $\hat{\theta}$. A lexicon $L$ is $n$-th order if $\text{ord}(L) = \max\{\text{ord}(\sigma(p)) \mid p \in \varnothing_0\} \leq n$.

**Definition 13** An abstract categorial grammar (ACG) is a quadruple $G = \langle \Sigma_0, \Sigma_1, L, s \rangle$, where

- $\Sigma_0$ is a higher-order signature, called the abstract vocabulary,
- $\Sigma_1$ is a higher-order signature, called the object vocabulary,
- $L$ is a linear lexicon from $\Sigma_0$ to $\Sigma_1$,
- $s \in \varnothing_0$ is called the distinguished type.

We sometimes call the triple $\langle a, \tau_0(a), L(a) \rangle$ for $a \in \varnothing_0$ a lexical entry, and specify an ACG by giving the set of lexical entries and the distinguished type.

**Definition 14** An ACG $G = \langle \Sigma_0, \Sigma_1, L, s \rangle$ generates two languages, the abstract language $A(G)$ and the object language $O(G)$, defined as

$$A(G) = \{ M \mid M \in \Lambda^{\text{lin}}(\Sigma_0) \text{ is a closed } \beta\eta\text{-normal term of type } s \},$$

$$O(G) = \{ |L(M)|_{\beta\eta} \mid M \in A(G) \}.$$

The abstract language can be thought of as a set of abstract grammatical structures, and the object language is regarded as the set of concrete forms obtained from these abstract structures and the lexicon. Thus, we simply say the language generated by an ACG for its object language. The term abstract categorial languages (ACLs) means the object languages of ACGs.

Though de Groote’s original definition of an ACG requires the lexicon to be linear, this paper allows the lexicon to be non-linear. We call an ACG whose lexicon is affine affine ACG, and denote the class of affine ACGs by $G^{\text{aff}}$. We then distinguish affine ACGs whose lexicons are linear, i.e., original ACGs, by calling them linear ACGs and let $G^{\text{lin}}$ denote the class of linear ACGs. Note that the abstract language always consists of linear terms, though an ACG is not necessarily linear. For each $G^* \in \{G^{\text{lin}}, G^{\text{aff}}\}$, $G^*(m, n)$ denotes the subclass of ACGs belonging to $G^*$ such that the order of the abstract vocabulary is at most $m$ and the order of the lexicon is at most $n$. An ACG is $m$-th order if it belongs to $G^*(m, n)$ for some $n$.

**Example 1** Let $\varnothing = o \rightarrow o$ and $M + N$ be an abbreviation of $\lambda z. M(Nz)$ if the types of $M$ and $N$ are $\varnothing$. Let us consider the affine ACG $G = \langle \Sigma_0, \Sigma_1, L, s \rangle$
with the following lexical entries:

<table>
<thead>
<tr>
<th>x ∈ Ξ₀</th>
<th>τ₀(x)</th>
<th>L(x)</th>
</tr>
</thead>
<tbody>
<tr>
<td>G</td>
<td>n</td>
<td>λv.v/cat/</td>
</tr>
<tr>
<td>M</td>
<td>n</td>
<td>λv.v/mouse/</td>
</tr>
<tr>
<td>J</td>
<td>np</td>
<td>λy.y/John/P₁</td>
</tr>
<tr>
<td>R</td>
<td>np → s</td>
<td>λx.x(λu.v + v/runs/</td>
</tr>
<tr>
<td>E</td>
<td>np → s</td>
<td>λx₁.x₂.x₃(λu.v + v/eats/</td>
</tr>
<tr>
<td>A</td>
<td>n → np</td>
<td>λz,y.(</td>
</tr>
<tr>
<td>L</td>
<td>n → np</td>
<td>λz,y.(</td>
</tr>
</tbody>
</table>

where each /xxx/ is a constant of type $\mathcal{A}$, $P_i$ denotes $\lambda u^n u^2 u_i$, $L(n) = (\mathcal{A}^2 \rightarrow \mathcal{A}) \rightarrow \mathcal{A}$, $L(np) = (\mathcal{A} \rightarrow (\mathcal{A}^2 \rightarrow \mathcal{A}) \rightarrow \mathcal{A}) \rightarrow \mathcal{A}$, $L(s) = \mathcal{A}$.

The object language $O(\mathcal{A})$ consists of terms representing some English sentences such as *John runs, all mice run, all cats eat a mouse*, and so on.

### 14.3 Linearization of Affine ACGs

While linear ACGs can generate languages consisting of linear terms only, affine ACGs can generate languages containing non-linear terms. Therefore, affine ACGs define a strictly richer class of languages than linear ACGs. However, since terms representing strings or trees are linear, affine terms in the object languages are not very interesting. This paper shows that for every $\mathcal{A} \in \mathbf{G}^{lin}(m, n)$, we can construct $\mathcal{A}' \in \mathbf{G}^{lin}(m, \max[2, n])$ such that

$$O(\mathcal{A}') = \{ P \in O(\mathcal{A}) \mid P \text{ is linear} \} \quad (14.8)$$

Moreover, in case of $m = 2$, we can find $\mathcal{A}' \in \mathbf{G}^{lin}(2, n)$ satisfying the equation (14.8). Therefore extending the definition of an ACG to allow lexical entries affine does not enrich the expressive power of ACGs in an essential way. Before proceeding with our construction, we mention a partially stronger result on the special case of this problem on string-generating second-order ACGs, obtained from Salvati’s work (Salvati, 2006). He presents an algorithm that converts a linear ACG $\mathcal{A} \in \mathbf{G}^{lin}(2, n)$ generating a string language into an equivalent LCFRS (via a deterministic tree-walking transducer). Even if an input is an affine ACG $\mathcal{A} \in \mathbf{G}^{all}(2, n)$, his algorithm still outputs an equivalent LCFRS. Since every LCFRS is encodable by a linear ACG belonging to $\mathbf{G}^{lin}(2, 4)$ (de Groote and Pogodalla, 2003, 2004), therefore this entails the following corollary.

---

3 A string $a_1 \ldots a_n$ on an alphabet $V$ is represented by $\lambda z. a_1(\ldots (a_n z) \ldots)$ $\in \mathbf{A}^{lin}(\Sigma_V)$ where $\Sigma_V = \langle \{a\}_V, \tau_V \rangle$ with $\tau_V(a) = \mathcal{A}$ for all $a \in V$ as in Example 1. Trees are constructed on some ranked alphabet. A ranked alphabet $(F, \rho)$, where $F$ is an alphabet and $\rho$ is a rank assignment on $F$, can be identified with a higher-order signature $\Sigma_{F, \rho} = \langle \{a\}_F, F, \tau_F \rangle$ such that $\tau_F(a) = \mathcal{A} \rightarrow \rho(a)$ if $\rho(a) = k$ for all $a \in F$, and a tree is identified with a variable-free (thus linear) term of the atomic type $\mathcal{A}$ in the obvious way.
Corollary 29 For every string-generating affine ACG \( \mathcal{G} \in G^{\mathbb{Z}}(2, n) \), there is a linear ACG \( \mathcal{G}' \in G^{\text{lin}}(2, 4) \) such that \( \mathcal{O}(\mathcal{G}) = \mathcal{O}(\mathcal{G}') \).

### 14.3.1 Basic Idea

We explain our basic idea for the linearization method for affine ACGs through a small example. Let us consider the affine ACG \( \mathcal{G} \) consisting of the following lexical entries:

<table>
<thead>
<tr>
<th>( x \in \mathcal{L}_0 )</th>
<th>( \tau_0(x) )</th>
<th>( \mathcal{L}(x) )</th>
</tr>
</thead>
<tbody>
<tr>
<td>A</td>
<td>( p \rightarrow s )</td>
<td>( \lambda w^{0 \rightarrow o}.w\alpha b^o )</td>
</tr>
<tr>
<td>B</td>
<td>( p )</td>
<td>( \lambda x^y . x )</td>
</tr>
</tbody>
</table>

where \( \mathcal{L}(s) = o \) and \( \mathcal{L}(p) = o^2 \rightarrow o \). Corresponding to \( AB \in \mathcal{A}(\mathcal{G}) \), we have \( a \in \mathcal{O}(\mathcal{G}) \) by

\[
\mathcal{L}(AB) \equiv (\lambda w^{0 \rightarrow o}.w\alpha b^o)(\lambda x^y . x) \rightarrow_{\beta} (\lambda x^y . x) \alpha^o \rightarrow_{\beta} \alpha^o .
\]

(14.9)

The occurrences of vacuous \( \lambda \)-abstraction \( \lambda y^o \) causes the deletion of \( b \) in (14.9). Such deleting operation is what we want to eliminate in order to linearize the affine ACG \( \mathcal{G} \). Let us retype \( \lambda y^o \) with \( \lambda \overline{y}^o \) and replace \( b^o \) with \( \overline{b}^o \) to indicate that they should be eliminated. Then (14.9) is decorated by bars as

\[
(\lambda w^{0 \rightarrow o}.w\overline{a} \overline{b}^o)(\lambda x^y . x) \rightarrow_{\beta} (\lambda x^y . x) \overline{a}^o \rightarrow_{\beta} \overline{a}^o .
\]

(14.10)

where we retype \( w^{0 \rightarrow o} \) with \( w^{0 \rightarrow \overline{a}} \), so that the whole term is well-typed.

In our setting, when a term has a barred type, it means that the term should be erased during \( \beta \)-reduction steps, and vice versa. By eliminating those barred terms and types from (14.10), we get

\[
(\lambda w^{0 \rightarrow o}.w\overline{a}^o)(\lambda x^o . x) \rightarrow_{\beta} (\lambda x^o . x) \overline{a}^o \rightarrow_{\beta} \overline{a}^o ,
\]

(14.11)

which solely consists of linear terms. Hence, the linear ACG \( \mathcal{G}' \) with the following lexical entries generates the same language as the original ACG \( \mathcal{G} \).

<table>
<thead>
<tr>
<th>( x \in \mathcal{L}_0 )</th>
<th>( \tau'_0(x) )</th>
<th>( \mathcal{L}'(x) )</th>
</tr>
</thead>
<tbody>
<tr>
<td>A</td>
<td>([p, o \rightarrow \overline{s} \rightarrow o] \rightarrow [s, o] )</td>
<td>( \lambda w^{0 \rightarrow o}.w\alpha )</td>
</tr>
<tr>
<td>B</td>
<td>([p, o \rightarrow \overline{s} \rightarrow o] )</td>
<td>( \lambda x^o . x )</td>
</tr>
</tbody>
</table>

where \([p, o \rightarrow \overline{s} \rightarrow o]\) and \([s, o]\) are new atomic types that are mapped to \( o \rightarrow o \) and \( o \) respectively, and \([s, o]\) is the distinguished type. We have \( \mathcal{L}(AB) = \mathcal{L}'(A'B') \). The term \( \lambda w^{0 \rightarrow o}.w\overline{a} \overline{b}^o \), which is led to \( \mathcal{L}'(A') \), is just one possible bar-decoration for \( \mathcal{L}(A) \). For instance, \( \lambda \overline{w}^{0 \rightarrow o}.w\overline{a} \overline{b}^o \) and \( \lambda w^{0 \rightarrow o}.w\overline{a} \overline{b}^o \) are also possible. Bars appearing in \( \lambda w^{0 \rightarrow o}.w\overline{a} \overline{b}^o \) predict that the sub-term \( \overline{a} \) will be erased, and \( \lambda w^{0 \rightarrow o}.w\overline{a} \overline{b}^o \) predicts that no sub-term of it will disappear. Our linearization method also produces lexical entries corresponding to those bar-decorations.
14.3.2 Formal Definition

We first give a formal definition of the set of possible bar-decorations on a type and a term. Hereafter, we fix a given affine ACG $G = (\Sigma_0, \Sigma_1, \mathcal{L}, s)$. Define $\Sigma_1 = (\mathcal{A}_1, \mathcal{P}_1, \mathcal{T}_1)$ by

$$\mathcal{A}_1 = \{ a \mid p \in \mathcal{A}_1 \}, \quad \mathcal{P}_1 = \{ c \mid \mathcal{C} \in \mathcal{G}_1 \}, \quad \mathcal{T}_1 = \{ \mathcal{T} \mid \mathcal{T}_1(\mathcal{C}) \mid \mathcal{C} \in \mathcal{G}_1 \},$$

where $\alpha \rightarrow \beta = \mathcal{T} \rightarrow \mathcal{P}$. Let $\Sigma_1' = (\mathcal{A}_1', \mathcal{P}_1', \mathcal{T}_1') = (\mathcal{A}_1 \cup \mathcal{A}_1', \mathcal{P}_1 \cup \mathcal{P}_1', \mathcal{T}_1 \cup \mathcal{T}_1').$

Here, we have the simple lexicon $\tilde{\mathcal{T}}$ from $\Sigma_1'$ to $\Sigma_1$ defined as

$$\tilde{\mathcal{T}} = p \mapsto p \text{ for } p \in \mathcal{P}_1, \text{ and } \tilde{\mathcal{C}} \equiv \tilde{\mathcal{C}} \equiv \mathcal{C} \text{ for } \mathcal{C} \in \mathcal{G}_1.$$

The set $\tilde{T}(\mathcal{A}_1')$ of possible bar-decorations on types is defined by

$$\tilde{T}(\mathcal{A}_1') = \{ \alpha \in T(\mathcal{A}_1') \mid \text{if } \beta_1 \rightarrow \cdots \rightarrow \beta_n \rightarrow \tilde{p} \text{ is a subtype of } \alpha \}$$

for some $\tilde{p} \in \tilde{\mathcal{T}}$, then $\beta_1, \ldots, \beta_n \in T(\mathcal{A}_1')$.

Actually, terms in $\Lambda^{aff}(\Sigma_1')$ that we are concerned with have types in $\tilde{T}(\mathcal{A}_1')$. The reason why we ignore types in $T(\mathcal{A}_1') - \tilde{T}(\mathcal{A}_1')$ is that if a term is bound to be erased, then so is every sub-term of it. For instance, if a variable $x$ has type $\alpha \rightarrow \beta = \mathcal{T} \rightarrow \mathcal{P}$, then the term $x^\alpha \rightarrow y^\beta$ has type $\mathcal{T}$, which, in our setting, means that it should disappear. But if $x^\alpha \rightarrow y^\beta$ disappears, so does $y^\beta$, which, therefore, should have type $\mathcal{T}$ to be consistent with our definition.

The set $\tilde{\Lambda}^{aff}(\Sigma_1')$ of possible bar-decorations on terms is the subset of $\Lambda^{aff}(\Sigma_1')$ such that $Q \in \Lambda^{aff}(\Sigma_1')$ iff

- every variable appearing in $Q$ has a type in $\tilde{T}(\mathcal{A}_1')$, and
- if $x^\alpha Q'$ is a sub-term of $Q$ and $x^\alpha$ does not occur free in $Q'$, then $\alpha \in T(\mathcal{A}_1')$.

We are not concerned with terms in $\Lambda^{aff}(\Sigma_1') - \tilde{\Lambda}^{aff}(\Sigma_1')$.

The following properties are easily seen:

- If $Q \in \tilde{\Lambda}^{aff}(\Sigma_1')$, then $T'(Q) \in \tilde{T}(\mathcal{A}_1')$.
- If $T'(Q) \in T(\Sigma_1')$ for $Q \in \tilde{\Lambda}^{aff}(\Sigma_1')$, every sub-term of $Q$ is in $\Lambda^{aff}(\Sigma_1')$.
- If $Q \in \tilde{\Lambda}^{aff}(\Sigma_1')$ and $Q \rightarrow_\beta Q'$, then $Q' \in \tilde{\Lambda}^{aff}(\Sigma_1')$.

For each $x \in T(\mathcal{A}_1')$ and $P \in \Lambda^{aff}(\Sigma_1')$, $\Phi$ gives the set of possible bar-decorations on them:

$$\Phi(x) = \{ \beta \in T(\mathcal{A}_1') \mid \tilde{\beta} = x \},$$

$$\Phi(P) = \{ Q \in \tilde{\Lambda}^{aff}(\Sigma_1') \mid \tilde{Q} = P \}.$$

In other words, $\Phi$ and $\tilde{\mathcal{T}}$ are inverse of each other, if we disregard types in $T(\mathcal{A}_1') - \tilde{T}(\mathcal{A}_1')$ and terms in $\Lambda^{aff}(\Sigma_1') - \tilde{\Lambda}^{aff}(\Sigma_1')$. 
Secondly, we eliminate barred subtypes from \( \alpha \in \overline{T(\mathcal{A}_1)} - T(\mathcal{A}_1) \) and barred sub-terms from \( Q \in \overline{\Lambda}^{\text{aff}}(\Sigma_1) - \Lambda^{\text{aff}}(\Sigma_1) \). Let us define \((\alpha)^\dagger\) and \((Q)^\dagger\) as follows:

\[
(\alpha \to \beta)^\dagger = \begin{cases} 
(\alpha)^\dagger \to (\beta)^\dagger & \text{if } \alpha \notin T(\mathcal{A}_1), \\
(\beta)^\dagger & \text{if } \alpha \in T(\mathcal{A}_1), 
\end{cases}
\]

\[
(\lambda \alpha.Q)^\dagger = \begin{cases} 
\lambda(\alpha)^\dagger.(Q)^\dagger & \text{if } \alpha \notin T(\mathcal{A}_1), \\
(Q)^\dagger & \text{if } \alpha \in T(\mathcal{A}_1), 
\end{cases}
\]

\[
(Q_1.Q_2)^\dagger = \begin{cases} 
(Q_1)^\dagger(Q_2)^\dagger & \text{if } \tau_1'(Q_2) \notin T(\mathcal{A}_1), \\
(Q_1)^\dagger & \text{if } \tau_1'(Q_2) \in T(\mathcal{A}_1). 
\end{cases}
\]

The following properties are easily seen \((\alpha)^\dagger \in T(\mathcal{A}_1)\) and \((Q)^\dagger \in \overline{\Lambda}^{\text{aff}}(\Sigma_1)\):

- \((\alpha)^\dagger \in T(\mathcal{A}_1)\) and \((Q)^\dagger \in \Lambda^{\text{lin}}(\Sigma_1)\),
- \(\tau_1((Q)^\dagger) = (\tau_1'(Q))^\dagger\),
- If \(Q\) is \(\beta\)-normal, then so is \((Q)^\dagger\),
- \(Q =_\beta Q'\) implies \((Q)^\dagger =_\beta (Q')^\dagger\).

**Lemma 30** For every closed term \(Q \in \overline{\Lambda}^{\text{aff}}(\Sigma_1)\), \(\tau_1'(Q) \in T(\mathcal{A}_1)\) iff \((Q)^\dagger =_\beta Q =_\beta \overline{Q}.

**Lemma 31** For every closed term \(P \in \Lambda^{\text{aff}}(\Sigma_1)\), \([P]_\beta\) is linear iff there is \(Q \in \Phi(P)\) whose type is in \(T(\mathcal{A}_1)\).

**Second-Order Case**

We say that an abstract atomic type \(p \in \mathcal{A}_0\) is useless if there is no \(M \in \mathcal{A}(\mathcal{F})\) that has a sub-term whose type contains \(p\). An abstract constant \(a \in \mathcal{C}_0\) is useless if there is no \(M \in \mathcal{A}(\mathcal{F})\) containing \(a\). If an ACG is second-order, it is easy to check whether the abstract vocabulary contains useless atomic types or constants, and if so, we can eliminate useless abstract atomic types and constants. This can be done in a way similar to the elimination of useless nonterminal symbols and production rules from a context-free grammar.

**Definition 15** Let \(\mathcal{G} = \langle \Sigma_0, \Sigma_1, \mathcal{L}, s \rangle\) be a second-order ACG that has no useless abstract atomic types or constants. We define \(\mathcal{G}' = \langle \Sigma'_0, \Sigma_1, \mathcal{L}', [s, \mathcal{L}(s)] \rangle\)
as follows: define $\Sigma'_0 = (\mathcal{A}_0', \mathcal{C}_0', \tau'_0)$ by
\[
\mathcal{A}_0' = \{ [p, \beta] \mid p \in \mathcal{A}_0, \beta \in \Phi(\mathfrak{L}(p)) - T(\mathfrak{A}_1) \},
\mathcal{C}_0' = \{ \{a, Q\} \mid a \in \mathcal{C}_0, Q \in \Phi(\mathfrak{L}(a)) - \Lambda^\text{aff}(\mathcal{A}_0) \},
\tau'_0 = \{ \{a, Q\} \mapsto ([\tau_0(a), \tau'_1(Q)])^\dagger \},
\]
where $([p, \beta])^\dagger = [p, \beta]$,
\[
([\alpha \rightarrow \gamma, \beta \rightarrow \delta])^\dagger = \begin{cases}
([\alpha, \beta])^\dagger \rightarrow ([\gamma, \delta])^\dagger & \text{if } \beta \notin T(\mathfrak{A}_1), \\
([\gamma, \delta])^\dagger & \text{if } \beta \in T(\mathfrak{A}_1),
\end{cases}
\]
and $\mathfrak{L}'$ by
\[
\mathfrak{L}'([p, \beta]) = (\beta)^\dagger, \quad \mathfrak{L}'([a, Q]) = (Q)^\dagger.
\]
$\mathfrak{G}'$ is linear, but it may contain useless abstract atomic types or constants. The linearized ACG $\mathfrak{G}'$ for $\mathfrak{G}$ is the result of eliminating all the useless abstract atomic types and constants from $\mathfrak{G}'$.

**Lemma 32.** Let $\mathfrak{G}$ and $\mathfrak{G}'$ be as in Definition 15.

For every variable-free $M \in \Lambda^\text{lin}(\Sigma_0)$ of an atomic type and every $Q \in \Phi(\mathfrak{L}(M)) - \Lambda^\text{aff}(\mathcal{A}_0)$, there is $N \in \Lambda^\text{lin}(\Sigma_0)$ such that $\tau'_0(N) = [\tau_0(M), \tau'_1(Q)]$ and $\mathfrak{L}'(N) \equiv (Q)^\dagger$.

Conversely, for every variable-free $N \in \Lambda^\text{lin}(\Sigma_0')$ of an atomic type, there are $M \in \Lambda^\text{lin}(\Sigma_0)$ and $Q \in \Phi(\mathfrak{L}(M)) - \Lambda^\text{aff}(\mathcal{A}_0)$ such that $\tau'_0(N) = [\tau_0(M), \tau'_1(Q)]$ and $\mathfrak{L}'(N) \equiv (Q)^\dagger$.

**Theorem 33.** For every affine ACG $\mathfrak{G} \in \mathbf{G}^\text{aff}(2, n)$, there is a linear ACG $\mathfrak{G}' \in \mathbf{G}^\text{lin}(2, n)$ such that $O(\mathfrak{G}') = \{ P \in O(\mathfrak{G}) \mid P \text{ is linear} \}$.

**Proof.** Use Lemmas 31, 32, and 30.

De Groote and Pogodalla (2003, 2004) have presented encoding methods for linear CFTGs and LCFRSs by linear ACGs. Their methods can also be applied to non-duplicating CFTGs and MCFGs.

**Example 2.** Let a non-duplicating CFTG $G$ consist of the following productions:

\[
S \rightarrow P(a, b), \quad P(x_1, x_2) \rightarrow P(c(x_1), c(S)) \mid d(x_1, x_2),
\]
where the ranks of $S$, $P$, $a$, $b$, $c$, $d$ are 0, 2, 0, 0, 1, 2, respectively. De Groote and Pogodalla’s method transforms $G$ into the following affine ACG $\mathfrak{G}$:

<table>
<thead>
<tr>
<th>$x \in \mathcal{C}_0$</th>
<th>$\tau_0(x)$</th>
<th>$\mathfrak{L}(x)$</th>
</tr>
</thead>
<tbody>
<tr>
<td>$A$</td>
<td>$p \rightarrow s$</td>
<td>$\lambda y_0^{\epsilon_s \rightarrow \epsilon_0} y_2 a b^p$</td>
</tr>
<tr>
<td>$B$</td>
<td>$s \rightarrow p \rightarrow p$</td>
<td>$\lambda y_0^{\epsilon_s \rightarrow \epsilon_0} x_1 x_2 y_2 (c(x_1)^{\epsilon_0} x_1) (c(x_1)^{\epsilon_0} y_2)$</td>
</tr>
<tr>
<td>$C$</td>
<td>$p$</td>
<td>$\lambda x_0^{\epsilon_p \rightarrow \epsilon_0} x_2 d^{\epsilon_p \rightarrow \epsilon_0} x_1 x_2$</td>
</tr>
</tbody>
</table>

$^4$The notation adopted here follows de Groote and Pogodalla.
When we apply the linearization method given in Definition 15 to $\mathcal{G}$, we get the following linear ACG $\mathcal{G}'$ whose distinguished type is $[s, o]$:

<table>
<thead>
<tr>
<th>$x \in \mathcal{G}'_0$</th>
<th>$\mathcal{L}'(x)$</th>
</tr>
</thead>
<tbody>
<tr>
<td>$[\mathcal{A}, L_y^{o \mapsto a}, y, \lambda]$</td>
<td>$L_y^{o \mapsto a}$</td>
</tr>
<tr>
<td>$[\mathcal{B}, L_y^{o \mapsto a}, y, \lambda]$</td>
<td>$L_y^{o \mapsto a}$</td>
</tr>
<tr>
<td>$[\mathcal{C}, P, [C, P]]$</td>
<td>$P(C, x_1, x_2)$</td>
</tr>
</tbody>
</table>

The linearized ACG $\mathcal{G}'$ is actually the encoding of the linear CFTG $G'$ consisting of the following productions:

\[
S \rightarrow P(a, b) \mid P'(a), \quad P'(x_1) \rightarrow P(c(x_1), c(S)) \mid P'(c(x_1)),
\]
\[
P(x_1, x_2) \rightarrow d(x_1, x_2),
\]

where the ranks of nonterminals $S, P, P'$ are 0, 2, 1, respectively. $G, \mathcal{G}, \mathcal{G}'$, and $G'$ generate the same tree language.

The following corollary generalizes the result by Fujiyoshi (2005), which covers the monadic case only.

**Corollary 34** For every non-duplicating CFTG $G$, there is a linear CFTG $G'$ such that $G$ and $G'$ generate the same tree language.

Let $\mathcal{G}$ be the affine ACG that encodes an MCFG $G$. The linearized ACG $\mathcal{G}'$ is indeed in the form that is the encoding of an LCFRS\(^5\) (but $\mathcal{G}$ is not). Therefore, our result covers the following theorem shown by Seki et al. (1991).

**Corollary 35** For every MCFG $G$, there is an LCFRS $G'$ such that the languages generated by $G$ and $G'$ coincide.

**Third or Higher-Order Case**

Definition 15 itself does not depend on the order of the given affine ACG except that in the general case, we do not know how to find and eliminate useless abstract atomic types and constants. For the general case, however, the linearized ACG given in Definition 15 may generate a strictly larger language.

\(^5\)The LCFRS obtained from an MCFG through our linearization method may have nonterminals of rank 0. The reason why usual definitions of an LCFRS do not allow nonterminals to have rank 0 is just to avoid redundancy. Mathematically speaking, allowing or disallowing nonterminals of rank 0 does not matter at all.
than the original affine ACG. In the remainder of this paper, we present a linearization method for general affine ACGs.

**Example 3** Suppose that an affine ACG \( \mathcal{G} \in \mathcal{G}^{\text{aff}}(3, 1) \) consists of the following lexical entries:

<table>
<thead>
<tr>
<th>( x \in \mathcal{E}_0 )</th>
<th>( \tau_0(x) )</th>
<th>( \mathcal{L}(x) )</th>
</tr>
</thead>
<tbody>
<tr>
<td>A</td>
<td>( q )</td>
<td>#</td>
</tr>
<tr>
<td>B</td>
<td>( p \to q \to q )</td>
<td>( \lambda y^p z^o b^p z)</td>
</tr>
<tr>
<td>C</td>
<td>( q \to s )</td>
<td>( \lambda z^o )</td>
</tr>
<tr>
<td>D</td>
<td>( (p \to s) \to s )</td>
<td>( \lambda x^{m \to n} a^{m \to n}(x b^n) )</td>
</tr>
</tbody>
</table>

We see \( \Omega(\mathcal{G}) = \{ a(\ldots(a(b(\ldots(b\#)\ldots))\ldots) \mid n \geq 0 \}. \) The linear ACG \( \mathcal{G}' \) by Definition 15 consists of the following lexical entries:

<table>
<thead>
<tr>
<th>( x \in \mathcal{E}_0' )</th>
<th>( \tau_0'(x) )</th>
<th>( \mathcal{L}'(x) )</th>
</tr>
</thead>
<tbody>
<tr>
<td>( [B, \lambda y^p z^o b^p z] )</td>
<td>( [q, o] \to [q, o] )</td>
<td>( \lambda z^o )</td>
</tr>
<tr>
<td>( [C, \lambda z^o )</td>
<td>( [q, o] \to [s, o] )</td>
<td>( \lambda z^o )</td>
</tr>
<tr>
<td>( [D, \lambda x^{m \to n} \cdot a(x \theta)] )</td>
<td>( [s, o] \to [s, o] )</td>
<td>( \lambda x^o )</td>
</tr>
</tbody>
</table>

The last lexical entry is useless. We have

\[
\Omega(\mathcal{G}') = \{ a(\ldots(a(b(\ldots(b\#)\ldots))\ldots) \mid m, n \geq 0 \} = \Omega(\mathcal{G}).
\]

Though any term of type \( p \) that is the first argument of an occurrence of \( B \) is bound to be erased in the original ACG \( \mathcal{G} \), we cannot ignore the occurrence of the type \( p \), because that occurrence of \( p \) balances the numbers of occurrences of \( B \) and \( D \) in a term in \( \mathcal{A}(\mathcal{G}) \).

Our new linearization method gives the linear ACG \( \mathcal{G}'' \) consisting of the following lexical entries (useless lexical entries are suppressed):

<table>
<thead>
<tr>
<th>( x \in \mathcal{E}_0'' )</th>
<th>( \tau_0''(x) )</th>
<th>( \mathcal{L}''(x) )</th>
</tr>
</thead>
<tbody>
<tr>
<td>( [B, \lambda y^p z^o b^p z] )</td>
<td>( [q, o] )</td>
<td>#</td>
</tr>
<tr>
<td>( [B, \lambda y^p z^o b^p z] )</td>
<td>( [p, o] \to [q, o] \to [q, o] )</td>
<td>( \lambda y^{m \to n} \cdot z^o(bz) )</td>
</tr>
<tr>
<td>( [C, \lambda z^o )</td>
<td>( [q, o] \to [s, o] )</td>
<td>( \lambda z^o )</td>
</tr>
<tr>
<td>( [D, \lambda x^{m \to n} \cdot a(x \theta)] )</td>
<td>( [(p, o] \to [s, o]) \to [s, o] )</td>
<td>( \lambda x^{(m \to n) \cdot o} \cdot a(x \theta) )</td>
</tr>
</tbody>
</table>

where \([p, \theta]\) is mapped to \( o \to o \). We have \( \Omega(\mathcal{G}) = \Omega(\mathcal{G}'') \).

Now, we give the formal definition of our new linearization method for general affine ACGs. For simplicity, we assume that \( \mathcal{A} = \{ o \} \) here, but it is possible to lift this assumption. The new linearized ACG \( \mathcal{G}'' \) has the form
\(G'' = (S_0', S_1, L'', [s, L(s)])\), where \(S_0'' = (\alpha, \beta, r_0''\})\) is defined by

\[
\begin{align*}
\alpha'' &= \{ [p, \beta] \mid p \in \alpha, \beta \in \Phi(L(p)) \}, \\
\beta'' &= \{ [a, Q] \mid a \in \alpha, Q \in \Phi(L(a)) \}, \\
r_0'' &= \{ [a, Q] \mapsto [r_0(a), r_1'(Q)] \}
\end{align*}
\]

where \([\alpha \rightarrow \gamma, \beta \rightarrow \delta] = [\alpha, \beta] \rightarrow [\gamma, \delta]\).

Here we have two simple lexicons \(L_0 : S_0'' \rightarrow \Sigma_0\) and \(L_1 : S_0'' \rightarrow \Sigma_1'\):

\[L_0([p, \beta]) = p, \quad L_0([a, Q]) = a, \quad L_1([p, \beta]) = \beta, \quad L_1([a, Q]) = Q.\]

We have \(L_1(N) \equiv L \circ L_0(N)\) for \(N \in \Lambda^{lin}(S_0'')\). For any \(M \in \Lambda^{lin}(\Sigma_0)\) and \(Q \in \Phi(L(M))\), one can find a term \(\chi(M, Q) \in \Lambda^{lin}(S_0'')\) such that \(L_0(\chi(M, Q)) \equiv M\) and \(L_1(\chi(M, Q)) \equiv Q\).

**Lemma 36** For every \(Q \in \Lambda^{al}(\Sigma_1)\) and \(\alpha \in T(\alpha_0)\), the following statements are equivalent:

1. There is \(M \in \Lambda^{lin}(\Sigma_0)\) of type \(\alpha\) such that \(L(M) \equiv \tilde{Q}\).
2. There is \(N \in \Lambda^{lin}(\Sigma_0'')\) of type \([\alpha, r_1'(Q)]\) such that \(L_1(N) \equiv Q\).

Lemmas 31 and 36 imply

\[(M \in A(\mathcal{G}) \upharpoonright \| L(M) \|_{\tilde{\theta}} \text{ is linear}) = \{ L_0(N) \mid N \in A(\mathcal{G}'') \}.\]

Since \((L_1(N))^\dagger = \tilde{\beta} L_1(N) \equiv L \circ L_0(N)\) for every \(N \in A(\mathcal{G}'')\) by Lemma 30, it is enough to define a new lexicon \(L''\) so that

\[
L''(N) \equiv_{\tilde{\beta_0}} (L_1(N))^\dagger \tag{14.12}
\]

for every \(N \in A(\mathcal{G}'')\).

We define the type substitution \(\sigma : \alpha'' \rightarrow T(\{o\})\) of \(L'' = (\sigma, \theta)\) as

\[
\sigma([p, \beta]) = \begin{cases} (\beta)^\dagger & \text{if } \beta \notin T(\mathcal{G}), \\
o & \text{if } \beta \in T(\mathcal{G}), \end{cases}
\]

Here we identify \(\sigma\) with its homomorphic extension. As a preparation for defining the term substitution \(\theta\) of \(L''\), we give three kinds of linear combinators. For \([\alpha, \beta] \in T(\alpha''\})\) such that \(\beta \in T(\mathcal{G})\), let \(\sigma([\alpha, \beta]) = \gamma_1 \rightarrow \cdots \gamma_m \rightarrow o \rightarrow o\) and \(\gamma_1 = \gamma_{1,1} \rightarrow \cdots \gamma_{1,k} \rightarrow o \rightarrow o\). \(Z^{\sigma([\alpha, \beta])}\) is a linear combinator of type \(\sigma([\alpha, \beta])\) defined as

\[
Z^{\sigma([\alpha, \beta])} \equiv A_{\gamma_1} \cdots A_{\gamma_m} \cdot R_1(R_2(\cdots (R_m \cdots )))
\]

where \(R_i \equiv y_i^{\gamma_i} \cdots y_m^{\gamma_m}\).

For each \([\alpha, \beta] \in T(\alpha''\})\) such that \(\beta \in T(\{o\}) \rightarrow T(\mathcal{G})\), we define two linear combinators \(X_{\alpha, \beta}'\) of type \(\sigma([\alpha, \beta]) \rightarrow (\beta)^\dagger\) and \(Y_{\alpha, \beta}'\) of type \((\beta)^\dagger \rightarrow \sigma([\alpha, \beta])\) by mutual induction. Let \([\alpha, \beta] = [\alpha_1, \beta_1] \rightarrow \cdots \rightarrow [\alpha_m, \beta_m] \rightarrow [p, \beta_0] \in \alpha''\) and the set \([1, \ldots, m]\) be partitioned into two subsets \(I\) and \(J\) so
that $\beta_i \notin \mathcal{T}(\overline{\mathcal{G}})$ iff $i \in I$. Let $I = \{i_1, \ldots, i_k\}$ $(i_j < i_{j+1})$ and $J = \{j_1, \ldots, j_l\}$.
Let
\[
X^{\beta_i}_n \equiv \lambda y^{\tau(x_{\alpha,\beta})(\beta_i)} X^{(\beta_i)}_{i_1} \cdots X^{(\beta_i)}_{i_k} y^{\tau(x_{\alpha,\beta})} P_1 \cdots P_m
\]
where $P_i \equiv \begin{cases} Y^{(\beta_i)}_{m}(\beta_i) & \text{if } i \in I, \\ Z^{\tau(\alpha,\beta)} & \text{if } i \in J, \end{cases}$ and
\[
\gamma^{\beta_i}_m \equiv \lambda x^{(\beta_i)}_{1} y^{\tau(x_{\alpha,\beta})(\beta_i)} \cdots y^{\tau(x_{\alpha,\beta})(\beta_i)} M_j \ldots (M_j(x^{(\beta_i)}_{L_1} \ldots L_n \bar{z}) \ldots )
\]
where $\bar{z}$ is short for $z^{(\beta_i)}_{1} \cdots z^{(\beta_i)}_{n}$ for $(\beta_0)^{\gamma} = \gamma_1 \rightarrow \cdots \rightarrow \gamma_n \rightarrow o$, and
\[
\begin{align*}
L_i & \equiv X^{\beta_i}_{i_1} y^{\tau(x_{\alpha,\beta})}_{i_1} & \text{for } i \in I, \\
M_i & \equiv Z^{\tau(\alpha,\beta)}_{i} & \text{for } i \in J.
\end{align*}
\]
Note that if $\alpha, \beta = [p, \beta_0] \in \mathcal{G}''_0$, then $X^{\beta_0}_p \equiv_{\beta_0} X^{\beta_0}_p \equiv_{\beta_0} \lambda x^{(\beta_0)} \gamma$.

Now, we give a new linearization method as follows.

**Definition 16** For a given affine ACG $\mathcal{G}$, we define a new linear ACG as $\mathcal{G}'' = (\Sigma''_0, \Sigma, \mathcal{L}'', [\cdot, \mathcal{L}'(\cdot)])$, where $\mathcal{L}'' = (\alpha, \theta)$ for $\alpha$ as above and
\[
\theta([a, Q]) \equiv \begin{cases} \{y^{\tau(s_m)(Q)}_{\gamma(a)} & \text{if } \tau(s_m)(Q) \notin \mathcal{T}(\overline{\mathcal{G}}), \\ Z^{\tau(s_m)(\gamma(a), Q)} & \text{if } \tau(s_m)(Q) \in \mathcal{T}(\overline{\mathcal{G}}). \end{cases}
\]

If $\mathcal{G} \in \mathcal{G}_0^{aff}(m, n)$, then $\mathcal{G}'' \in \mathcal{G}_0^{lin}(m, \max(2, n))$.

**Lemma 37** Given $N \in \Lambda(\Sigma''_0)$ of type $[\alpha, \beta]$ such that $\beta \notin \mathcal{T}(\overline{\mathcal{G}})$ and $\mathcal{L}'(N) \in \Lambda^{lin}(\Sigma)$, we have
\[
(\mathcal{L}'(N))^+ \equiv_{\beta_0} X^{\beta_0}_n \mathcal{L}''(N)\phi_N
\]
where $\phi_N$ is the substitution on the free variables of $\mathcal{L}''(N)$ such that
\[
x^{\tau(x_{\alpha,\beta})} \phi_N = \begin{cases} Y^{\tau(x_{\alpha,\beta})}_n & \text{if } x \text{ has the type } [\alpha, \beta] \text{ in } N \text{ and } \beta \notin \mathcal{T}(\overline{\mathcal{G}}), \\ Z^{\tau(x_{\alpha,\beta})} & \text{otherwise.} \end{cases}
\]

**Theorem 38** For every affine ACG $\mathcal{G} \in \mathcal{G}_0^{aff}(m, n)$, there is a linear ACG $\mathcal{G}'' \in \mathcal{G}_0^{lin}(m, \max(2, n))$ such that $\mathcal{O}(\mathcal{G}'') = \{P \in \mathcal{O}(\mathcal{G}) \mid P \text{ is linear}\}$.

**Proof.** Lemma 37 entails the equation (14.12). \qed

### 14.4 Concluding Remarks

We have shown that the generative capacity of linear ACGs is as rich as that of affine ACGs, that is, the non-deletion constraint on linear ACGs is superficial. Our linearization method, however, increases the size of the given grammar exponentially due to the definition of $\Phi$, so there may still exist an advantage of allowing deleting operations in the ACG formalism. For instance, the
atomic type \( np \) of the abstract vocabulary of the ACG in Example 1 will be divided up into three new atomic types which correspond to noun phrases as third person singular subjects, plural subjects, and objects, respectively.

One attractive feature of ACGs is that they can be thought of as a generalization of several well-established grammar formalisms (de Groote, 2002, de Groote and Pogodalla, 2003, 2004). This paper demonstrates that the ACG formalism also generalizes some “operation” on those grammars, namely, conversion from non-duplicating grammars into non-duplicating and non-deleting ones.

Recall that Fisher (1968a,b) showed that every CFTG has a corresponding non-deleting CFTG whose string IO-language is equivalent. As a generalization of his result, the author conjectures that one can eliminate vacuous \( \lambda \)-abstraction from semi-affine ACGs preserving the orders of the abstract vocabularies and the lexicons, where a term is semi-affine if for every free variable \( x \) of any sub-term, either \( x \) occurs at most once, or \( x \) has at most second-order type. Actually, every CFTG has a corresponding semi-affine ACG such that the tree IO-language of the CFTG coincides with the object language of the ACG, and the semi-affine ACG encoding a non-deleting CFTG has no vacuous \( \lambda \)-abstraction. If the conjecture is correct, this implies that every CFTG has a corresponding non-deleting CFTG whose tree/string IO-language is equivalent.

Acknowledgment

The author is grateful to Makoto Kanazawa for initiating this research and giving advice throughout this work. The author would like to thank to Sylvain Salvati for his invaluable comments on the draft of this paper. In particular, he inspired the author to get the conjecture stated in the last section.

References


List of contributors

Miguel A. Alonso
Departamento de Computación, Facultad de Informática
Universidade da Coruña, Spain

Maxime Amblard
LaBRI, Université Bordeaux I
Bordeaux, France

Houda Anoun
LaBRI, Université Bordeaux I
Bordeaux, France

John Blatz
Department of Computer Science, Johns Hopkins University
Baltimore, MD, USA

Maria Bulińska
Faculty of Mathematics and Computer Science, Uniwersytet Warmińsko-Mazurski
Olsztyn, Poland

Jason Eisner
Department of Computer Science, Johns Hopkins University
Baltimore, MD, USA

Josef van Genabith
National Centre for Language Technology, School of Computing, Dublin City University
Dublin, Ireland

Carlos Gómez-Rodríguez
Departamento de Computación, Facultad de Informática
Universidade da Coruña, Spain

Laura Kallmeyer
Eberhard-Karls Universität Tübingen
Tübingen, Germany

**Stephan Kepser**
Eberhard-Karls Universität Tübingen
Tübingen, Germany

**Aleksandra Kislak-Malinowska**
Uniwersytet Warmińsko-Mazurski
Olsztyn, Poland

**Alain Lecomte**
Université Pierre Mendes-France (Grenoble II)
Grenoble, France

**Rebecca Nesson**
Division of Engineering and Applied Sciences, Harvard University
Cambridge, MA, USA

**Sylvain Salvati**
National Institute of Informatics
Tokyo, Japan

**Stuart Shieber**
Division of Engineering and Applied Sciences, Harvard University
Cambridge, MA, USA

**Edward P. Stabler**
Linguistics Department, UCLA
Los Angeles, CA, USA

**Jesse Tseng**
CNRS/Loria
Vandoeuvre-lès-Nancy, France

**Manuel Vilares**
Departamento de Computación, Facultad de Informática
Universidade da Coruña, Spain

**Ryo Yoshinaka**
National Institute of Informatics, University of Tokyo
Tokyo, Japan