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# Probabilistic Constraint Logic Programming

Stefan Riezler

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# Probabilistic Constraint Logic Programming

Formal Foundations of Quantitative and Statistical Inference in Constraint-Based Natural Language Processing

Stefan Riezler

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Lehrstuhl für Theoretische Computerlinguistik Institut für Maschinelle Sprachverarbeitung Universität Stuttgart Azenbergstraße 12 70174 Stuttgart

www: http://www.ims.uni-stuttgart.de/projekte/gramotron/ email: gramotron@ims.uni-stuttgart.de

# Probabilistic Constraint Logic Programming Formal Foundations of Quantitative and Statistical Inference in Constraint-Based Natural Language Processing

von

Stefan Riezler

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Mitberichterstatter: PhD Steven Abney, AT&T Labs, Florham Park, NJ
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#### Abstract

Structural ambiguity in linguistic analyses is a severe problem for natural language processing. For grammars describing a nontrivial fragment of natural language, every input of reasonable length may receive a large number of analyses, many of which are implausible or spurious. This problem is even harder for highly complex constraint-based grammars. Whereas the mathematical foundation of such grammars as instances of constraint logic programming is clear enough, there is so far no mathematically well-defined method for extending constraint logic programs by using weights to introduce graded distictions between analyses. Previous approaches to ambiguity resolution for context-dependent natural language processing models either are tailored to specific applications and based on uncertain mathematical grounds, or they are sufficiently well-defined and expressive but infeasible in practice.

In this thesis, we present two approaches to a rigorous mathematical and algorithmic foundation of quantitative and statistical inference in constraint-based natural language processing. The first approach, called quantitative constraint logic programming, is conceptualized in a clear logical framework, and presents a sound and complete system of quantitative inference for definite clauses annotated with subjective weights. This approach combines a rigorous formal semantics for quantitative inference based on subjective weights with efficient weight-based pruning for constraint-based systems. The second approach, called probabilistic constraint logic programming, introduces a log-linear probability distribution on the proof trees of a constraint logic program and an algorithm for statistical inference of the parameters and properties of such probability models from incomplete, i.e., unparsed data. The possibility of defining arbitrary properties of proof trees as properties of the log-linear probability model and efficiently estimating appropriate parameter values for them permits the probabilistic modeling of arbitrary context-dependencies in constraint logic programs. The usefulness of these ideas is evaluated empirically in a small-scale experiment on finding the correct parses of a constraint-based grammar. In addition, we address the problem of computational intractability of the calculation of expectations in the inference task and present various techniques to approximately solve this task. Moreover, we present an approximate heuristic technique for searching for the most probable analysis in probabilistic constraint logic programs.

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# Chapter 1

# Introduction

This thesis presents a novel mathematical treatment of the problem of structural ambiguity in constraint-based natural language processing (NLP). This problem will be attacked from two different angles. On the one side we will present a novel formalism for quantitative constraint-based inference with subjective weights. On the other side we will approach this problem by novel methods for statistical inference and probabilistic modeling for constraint-based NLP.

In this chapter we introduce the general problem of structural ambiguity and a general solution to this problem, namely weighted grammars. Furthermore, we will specify the notion of constraint-based NLP and sketch the general idea of the two different approaches to ambiguity resolution for weighted constraint-based grammars which constitute the main contribution of this thesis.

## 1.1 Overview

Following this introduction, Chap. 2 discusses the formal framework in which the informal notion of constraint-based NLP will be dealt with in the course of this thesis. To this end, we discuss the formal basics of Constraint Logic Programming (CLP), which is used here to provide an operational treatment of various declarative constraint-based grammars. This is done by an embedding of the logical description languages of such grammars into a CLP scheme, yielding Constraint Logic Grammars (CLGs).

Chap. 3 presents a quantitative extension of CLP which allows us to assign subjective numerical weights to the structural components of a constraint logic program. We present a sound and complete system for quantitative inference with such subjective weights based on concepts of fuzzy set algebra. Furthermore, the general concepts of quantitative CLP will be exemplified with a simple quantitative CLG and we will show how the search technique of alpha-beta pruning can be adapted to efficiently finding the best parse in quantitative CLGs. A completely different approach to weighted CLP is presented in Chap. 4. Here, instead of concentrating on a formal specification of the handling of subjective weights, the aim is to use methods of probabilistic modeling and statistical inference to automatically induce weights from empirical data. We introduce a powerful log-linear probability model for CLP and present a novel technique for statistical inference of the parameters and properties of such models from incomplete training data. We show monotonicity and convergence of the algorithm to the desired maximum likelihood estimates and discuss various methods for approximate computation for the inference task. We present an instantiation of probabilistic CLP to a simple probabilistic CLG and show how the structure of the probabilistic model can be used to guide the search for the most probable analysis. Furthermore, the main concepts of this statistical approach are evaluated empircally in a small experiment on finding the correct parses of a constraint-based grammar.

Chaps. 3 and 4, presenting the two different approaches to weighted CLP and CLGs, are conceptualized completely independent of each other. Whereas Chap. 3 is based upon the general concepts of Chap. 2, namely classical CLP with CLGs as a special instance, the work of Chap. 4 is entirely self-contained and even more general. That is, the presented methods of probabilistic modeling, statistical inference and approximate computation can easily be abstracted away from the CLP application to more general data structures.

Chap. 5 presents a summary of the work of this thesis, and compares the advantages and shortcomings of the two presented approaches relative to each other and relative to other approaches. Furthermore, directions of future work are sketched.

The rest of this chapter presents a motivation of the why and how of the work of this thesis.

## 1.2 A Practical Problem: Structural Ambiguity

Structural ambiguity is a practical problem for every grammar describing a nontrivial fragment of natural language. That is, for such grammars every input of reasonable length may receive a large number of different analyses, many of which are not in accord with human perceptions. The problem to be addressed is how to differentiate between these analysis and how to efficiently find the correct analysis out of the set of all possible ones.

A simple example illustrating the ubiquity and severity of the problem of structural ambiguity has been presented by Church and Patil (1982). Consider the following sentence with two PPs. It has the following two analyses in terms of PP-attachment:

(1) a. Put the block [in the box on the table].

b. Put [the block in the box] on the table.

If we have three PPs, the number of analyses is five.

- (2) a. Put the block [[in the box on the table] in the kitchen].
  - b. Put the block [in the box [on the table in the kitchen]].
  - c. Put [[the block in the box] on the table] in the kitchen.
  - d. Put [the block [in the box on the table]] in the kitchen.
  - e. Put [the block in the box] [on the table in the kitchen].

Continuing this list further, a number of more than thousand analyses is achieved quickly with only eight PPs. The pattern behind this list can be explained as a combinatorial growth of ambiguity in the number of PPs. This growth pattern follows the combinatorial principle of the the Catalan numbers, where Cat(n) describes the number of ways to parenthesize a sentence of length n, or equivalently the set of binary trees that can be constructed over nterminal elements<sup>1</sup>. Clearly, this pattern can be found also in other linguistic combinations such as conjuncts, nominal modifications, or relative clauses.

Whereas ambiguities of this kind are only problematic if the number of linguistic elements to be combined is large, there is another source of ambiguity depending simply on the number of analyses the grammar can produce at all. Let us consider the standard linguistic example sentence John saw Mary and the two analyses given below.

- (3) a.  $[\operatorname{John}_N [\operatorname{saw}_V \operatorname{Mary}_N]_{VP}]_S$ .
  - b.  $[[John_N \text{ saw}_N]_{NP} \text{ Mary}_N]_{NP}.$

Even if the first analysis is perfectly plausible and might be considered as the unique analysis of this sentence, the second analysis has to be accepted if the grammar also licenses other nominal modifications such as

(4) [[school<sub>N</sub> committee<sub>N</sub> ]<sub>NP</sub> meeting<sub>N</sub> ]<sub>NP</sub>.

Following Abney (1996), the second analysis furthermore can be given a perfectly plausible interpretation as the reference to a person named Mary who is associated with a kind of saw called John saw. Clearly, such spurious ambiguities may be characterized as resulting from rare usages of words and constructions, but they will appear in every grammar which covers a reasonable fragment of natural language and thus produces a large number of analyses. Furthermore, in most cases such spurious ambiguities cannot be given a plausible interpretation, but just have to be accepted as a side-effect of high coverage.

<sup>1</sup>The Catalan numbers are generated by the following formula:  $Cat(n) = \begin{pmatrix} 2n \\ n \end{pmatrix} - \begin{pmatrix} 2n \\ n-1 \end{pmatrix}$ .

Together combinatorial and spurious ambiguity can confront NLP systems with severe problems. Clearly, there is a need to distinguish more plausible analyses of an input form less plausible or even totally spurious ones. A practical and general approach to this problem is the use of weighted grammars for resolving structural ambiguities.

## **1.3** A Practical Solution: Weighted Grammars

We will approach the problem of structural ambiguity by using weighted grammars for ambiguity resolution. Weighted grammars can be characterized very generally as follows. They assign numerical values, called weights, to the structure-building components of a grammar and calculate the weight of an analysis from the weights of the structural components that make it up. The simple but effective assumption is to connect the plausibility of an analysis with its weight. That is, a ranking of analyses is defined by the weighted grammar, and more plausible analyses are differentiated from less plausible analyses in terms of their weights. The most plausible or correct analysis then is chosen from among the in-principle possible analyses by assuming the analysis with the greatest weight to be the correct one. Furthermore, when we are interested only in the highest weighted parse, the weight calculation scheme can be used to guide the search for the highest weighted parse efficiently instead of simply listing all possible parses and choosing the highest weighted one.

There are three basic problems to be solved for every weighted grammar to be a useful device in real-world NLP applicatons. These problems can be described by the following questions.

- 1. How can the values of the weights be obtained?
- 2. How should the weights be applied to the components of the grammar and how should the weight of an analysis be calculated from the weights of the components?
- 3. How can the structure of the weight calculation scheme be used to guide the search for the highest weighted analysis efficiently?

Clearly, the answers to these questions depend on each other and on the non-weighted framework to be extended. In the following we will sketch the basic ideas of two different approaches to answer these questions consistently for a framework of constraint-based systems.

# 1.4 Towards a Mathematical Foundation of Weighted Constraint-Based Grammars

The NLP systems of choice in this thesis are constraint-based grammars. The term constraintbased is a collective name for highly expressive frameworks for declarative description of natural language in terms of logical description languages. Throughout this thesis, the informal concept of constraint-based grammars will be replaced by the formal concept of constraint logic grammars. That is, constraint-based grammars are formalized here by an embedding of the logical description languages of such grammars into a CLP scheme, yielding CLGs as special applications of CLP. The advantages of this approach are on the one hand the (Turing-)power of the underlying logic, which is conceived as a welcome property to overcome the inadequacy of regular and context-free grammars for the description of natural language. On the other hand this approach permits an operational treatment of, e.g., the parsing problem for arbitrary constraint-based grammars in a consistent and unique way. Since CLGs can be seen as special applications of CLP, the mathematical work of this thesis will be based upon CLP in general, and CLGs will serve as running example illustrating the applicability of the general work to NLP. The reference to the general framework of CLP will generalize the results of this thesis in a welcome manner.

However, most CLP applications require some form of graded distinctions which are not provided by a classical CLP scheme. A very important example for this demand for gradedness is the task of structural ambiguity resolution in CLGs. A crucial assumption in this thesis is the claim that a framework of weighted CLP is the solution of choice for the ambiguity resolution problem for CLGs. In the following chapters we will present a rigorous mathematical formulation of two different approaches to weighted CLP and weighted CLGs.

The first approach we will present is motivated by the aim to give the grammar designer and implementer maximal freedom in choosing appropriate values for the weights of the weighted grammar. That is, the values of the weights are only restricted to be some quantities lying in a certain interval of real numbers. Such weights can be restricted to meet the axioms of probability theory, but there is no need to do so. Besides subjective probabilities, such quantities could be subjective preference values, or values obtained from experiments on preferences in human language processing, or values describing human judgements on degrees of grammticality, or others. In order to stress the generality of this approach to weighted CLP and weighted CLGs, we will henceforth refer to it as *quantitative CLP* and *quantitative CLGs*, respectively. The main task of this approach is to specify the questions of how to establish a proper weight calculation scheme for given values and of how to use such a scheme for efficient disambiguation. Since it is the grammar designer and implementer who has to specify the grammar and the weights, it makes sense to tie these two tasks together as closely as possible. That means, in the same way as the inference system of classical CLP is coupled with a clear formal semantics, one would like to relate a quantitative inference system to a quantitative formal semantics, instead of adding an extralogical calculation scheme to the well-defined logic of CLP. Thus the task to be addressed is to provide a precise, but yet simple formal semantics for quantitative inference in CLP. To this end, we present a formal semantics for quantitative CLP based upon the simple and intuitive concepts of fuzzy set algebra. This semantics and the corresponding sound and complete quantitative inference system furthermore are designed in a way which enables the search technique of alpha-beta pruning to be used quite directly for efficient disambiguation. Quantitative CLP then provides an efficient, well-defined quantitative deduction system, which can be adapted for specific applications by embedding specific constraint languages into CLP and attaching appropriate weights to them.

A completely different approach to weighted CLP and weighted CLGs is presented by our models of probabilistic CLP and probabilistic CLGs. The aim of this approach is to specify a probability distribution over the set of proof trees of CLP or the parses of CLGs, and to provide statistical methods to infer the values of the parameters of such probabilitic models from empirical data. For a given sample of training data and a parametric probability model, both the parameters of the probabilistic model and the properties of the model associated with these parameters can be induced automatically by methods of statistical inference. We present a highly expressive log-linear probability model for CLP, and a novel algorithm to infer the parameters and properties of log-linear models from incomplete data. We show monotonicity and convergence of the new algorithm and discuss methods for efficient approximate computation of the formulae involved in the algorithm. This algorithm is applicable to log-linear models in general, and especially provides the means for automatic and reusable training of arbitrary probabilistic constraint-based grammars from unparsed data. The usefulness of these concepts is shown empirically in a small-scale experiment on finding preferences in parse-data from a constraint-based grammar. Furthermore, we discuss the possibilities of using the structure of the probabilistic model to guide the search for the most probable proof tree or analysis, and present a heuristic search algorithm for this task. Clearly, in this setting a model-theoretic semantics for probabilistic inference is superfluous since the values of the probabilistic parameters are obtained by automatic statistical methods which are not manipulable by the user. Rather, we are interested in a stochastic semantics for CLP inference which is determined by the log-linear probability model together with the statistical methods for parameter estimation and property selection from given input data.

## 1.5 Bibliographical Note

Various parts of this thesis are based upon previously published work of the author. Chap. 3 is an extended version of Riezler (1996). Chap. 4 is based upon work presented in Riezler (1997), Riezler (1998a), Riezler (1998b), and Johnson, Geman, Canon, Chi, and Riezler (1999).

# Chapter 2

# Foundations: Basic Concepts of CLP and CLGs

In this chapter we report the central formal concepts of the CLP scheme of Höhfeld and Smolka (1988). In preparation for the following work we give some proofs missing in the original paper and present the CLP scheme in a slightly modified fashion. Furthermore, in order to prepare the running example of the next chapters, we report the main concepts of a feature-based constraint language for HPSG and show how to embed this constraint language into the CLP scheme, yielding feature-based CLGs.

### 2.1 Introduction and Overview

Constraint logic programming is a powerful extension of conventional logic programming (Lloyd 1987), and involves the incorporation of constraint languages and constraint solving methods into logic programming languages. The name CLP was first introduced by Jaffar and Lassez (1986) for a general framework of a logic programming language that is parametrized with respect to constraint language and a domain of computation, and yields soundness and completeness results for an operational semantics relying on a constraint solver for the employed constraint language. For example, conventional logic programming or Prolog is obtained from CLP by employing equations between first order terms as constraint language and by interpreting these equations in the Herbrand universe. In this case the operational semantics of SLD-resolution can be seen to rely on a constraint solver which solves term equations in the Herbrand universe. In this case the operational semantics of SLD-resolution can be seen to rely on a constraint solver which solves term equations in the Herbrand universe. In this case the operational semantics of SLD-resolution can be seen to rely on a constraint solver which solves term equations in the Herbrand universe. In this case the operational semantics of SLD-resolution can be seen to rely on a constraint solver which solves term equations in the Herbrand universe by term unification. Recent extensions, refinements, and various applications of CLP are discussed in Jaffar and Maher (1994). In the following we will rely on the general CLP scheme of Höhfeld and Smolka (1988), which has been shown to be a useful tool for our intended application of linguistic knowledge representation (see Dörre and Dorna

(1993), Götz (1995), Götz and Meurers (1995)).

The term constraint logic grammars expresses the connection between CLP and constraintbased grammars. That is, CLGs are understood as grammars formulated by means of a suitable logical language which can be used as a constraint language in the CLP scheme of Höhfeld and Smolka (1988). The idea behind this connection is to provide an operational treatment of purely declaratively specified grammars. This needs further explanation: Constraint-based grammars enable a clear model-theoretic characterization of linguistic objects by specifying grammars as sets of descriptions from a suitable logical description language, called the constraint language. The descriptions, called constraints, are stated as axioms required to be true of every object in the domain to be described, i.e., they constrain the admissible models of the grammar. The parsing problem (and similarly the generation problem) can be defined as follows: Given a set of axioms (encoding the grammar) and some constraint  $\phi$  (encoding the string/logical form we want to parse/generate from), we ask if there is some model of our axioms which satisfies  $\phi$ . Following Götz (to appear), we will call this the prediction problem.

A well-known subclass of these grammars widely used in computational linguistics are grammars based upon feature description languages such as simple PATR grammars (Shieber 1986) or more expressive grammars such as LFG (Bresnan and Kaplan 1982) or HPSG (Pollard and Sag 1994). Formalizations of the more or less informal notions of these grammars in terms of first-order languages were firstly presented by Smolka (1988) for PATR and by Johnson (1988) and King (1989), King (1994) for LFG and HPSG, respectively.

However, such model-theoretic approaches do not necessarily provide an operational interpretation of their declarative specifications. This may lead to problems with an operational treatment of model-theoretically well-defined problems such as parsing or generation. CLP provides one possible approach to an operational treatment of various such frameworks by embedding arbitrary logical languages into constraint logic programs. Definite clause specifications over such constraint languages then define grammars as constraint logic programs, i.e., as sets of axiomatic interpreted definite clauses. The prediction problem is in this setting as follows: Given a program  $\mathcal{P}$  (encoding a grammar) and a definite goal G (encoding the string/logical form we want to parse/generate from), we ask if we can infer an answer  $\varphi$  of G(which is a satisfiable constraint encoding an analysis) proving the implication  $\varphi \to G$  to be a logical consequence of  $\mathcal{P}$ .

For feature-based grammars an embedding of a logical language close to that of Smolka (1988) into the CLP scheme of Höhfeld and Smolka (1988) is done in the formalism CUF (Dörre and Eisele 1991; Dörre and Dorna 1993). This approach quite directly offers the operational properties of the CLP scheme, but unfortunately gives up the connection to the model-theoretic specifications of the underlying feature-based grammars. A different approach is given by Götz (1995), Götz and Meurers (1995), who defines an explicit translation from a logical language close to that of King (1994) into constraint logic programs. This trans-

lation procedure preserves the prediction problem by generating a constraint logic program  $\mathcal{P}(\mathcal{G})$  from a feature-based grammar  $\mathcal{G}$  in an explicit way. Other approaches to an operational semantics for the prediction problem of feature-based languages have been presented, e.g., by Carpenter (1992), Aït-Kaci, Podelski, and Goldstein (1993) or Götz (to appear). These approaches are tailored especially for specific feature-based languages and clearly suit the particular frameworks better than an embedding of the specific languages into a CLP scheme. However, under the CLP approach, arbitrary constraint-based grammars can receive an unique operational semantics by an embedding into definite clause specifications<sup>1</sup>.

We see the main advantage of the CLP approach in the possibility to rely on the wellunderstood paradigm of logic programming. This allows the resulting programs to run on existing architectures and to use well-known optimization techniques worked out in this area. The possibility to embed arbitrary constraint languages into the CLP scheme and the broad applicability of CLP itself should generalize the work of the following chapters in a welcome manner.

This chapter is organized as follows. In Sect. 2.2 we will report the main concepts of constraint logic programming following the CLP scheme of Höhfeld and Smolka (1988). As the work in the next chapters will build upon this scheme, we will reformulate the main definitions and propositions of Höhfeld and Smolka (1988) in a form convenient for the following discussions, and give some missing proofs which will be helpful to make this work parallel to the work of the next chapters.

In order to provide a concrete instantiation of this CLP scheme to constraint logic grammars, we will report in Sect. 2.3 a feature-based constraint language and show how this language can be embedded into the CLP scheme to yield feature-based CLGs.

## 2.2 Constraint Logic Programming

The scheme presented by Höhfeld and Smolka (1988) generalizes conventional logic programming (Lloyd 1987) and also the constraint logic programming scheme of Jaffar and Lassez (1986) to a scheme of definite clause specifications over arbitrary constraint languages. Relying on terminology well-known for conventional logic programming, Höhfeld and Smolka's generalization of the key result of conventional logic programming can be stated as follows: First, for every definite clause specification  $\mathcal{P}$  in the extension of an arbitrary constraint language  $\mathcal{L}$ , every interpretation of  $\mathcal{L}$  can be extended to a minimal model of  $\mathcal{P}$ . Second, the SLD-resolution method for conventional logic programming can be generalized to a sound and complete operational semantics for definite clause specifications, which are not restricted to

<sup>&</sup>lt;sup>1</sup>For example, an embedding of a the logical language for tree-description grammars of Rogers (1994) into the CLP scheme of Höhfeld and Smolka (1988) is given in Morawietz (1997).

Horn theories. In contrast to Jaffar and Lassez (1986), in this scheme constraint languages are not required to be sublanguages of first order predicate logic and do not have to be interpreted in a single fixed domain. Instead, a constraint is satisfiable if there is at least one interpretation in which it has a solution. This makes this scheme usable for a wider range of applications. Furthermore, such interpretations do not have to be solution compact<sup>2</sup>. This was necessary in Jaffar and Lassez (1986) to provide a sound and complete treatment of negation as failure. Höhfeld and Smolka (1988) do not include negation as failure but rather let the embedded constraint language provide for logical negation.

#### 2.2.1 Constraint Languages

A very general characterization of the concept of constraint language can be given as follows.

**Definition 2.1** ( $\mathcal{L}$ ). A constraint language  $\mathcal{L}$  consists of

- an  $\mathcal{L}$ -signature, specifying the non-logical elements of the alphabet of the language,
- a decidable infinite set VAR whose elements are called variables,
- a decidable set CON of  $\mathcal{L}$ -constraints which are pieces of syntax built from the  $\mathcal{L}$ -signature, the variables in VAR, and the logical elements of the alphabet of the language,
- a computable function V assigning to every constraint  $\phi \in \text{CON}$  a finite set  $V(\phi)$  of variables, the variables constrained by  $\phi$ ,
- a nonempty set of *L* -interpretations INT, where each *L* -interpretation *I* ∈ INT is defined w.r.t. a nonempty set *D*, the domain of *I*, and a set ASS of variable assignments α : VAR → *D*,
- a function [[·]]<sup>T</sup> mapping every constraint φ ∈ CON to a set [[φ]]<sup>T</sup> of variable assignments, the solutions of φ in I.
- Furthermore, a constraint  $\phi$  constraints only the variables in  $V(\phi)$ , i.e., if  $\alpha \in \llbracket \phi \rrbracket^{\mathcal{I}}$  and  $\beta$  is a variable assignment that agrees with  $\alpha$  on  $V(\phi)$ , then  $\beta \in \llbracket \phi \rrbracket^{\mathcal{I}}$ .

In order to state certain closure conditions on constraint languages, further definitions are necessary. The following definitions are made with respect to some given constraint language.

#### Definition 2.2.

• A renaming is a bijection VAR → VAR that is the identity except for finitely many exceptions.

 $<sup>^{2}</sup>$ That is, it is not necessary that every element of an interpretation must be obtainable as the unique solution of a possibly infinite set of constraints. See Jaffar and Lassez (1986).

- A constraint φ' is a ρ-variant of a constraint φ under a renaming ρ iff φ' = φρ, i.e., φ' is the constraint obtained from φ by simultaneously replacing each occurence of a variable X in φ by ρ(X) for all variables X in V(φ), and so [[φ]]<sup>I</sup> = [[φ']]<sup>I</sup><sub>αορ</sub> := {α ορ| α ∈ [[φ']]<sup>I</sup>, i.e., the function compositions of the solutions of φ' and a renaming ρ yield the solutions of φ, for all interpretations I.
- A constraint  $\phi'$  is a variant of a constraint  $\phi$  if there exists a renaming  $\rho$  s.t.  $\phi'$  is a  $\rho$ -variant of  $\phi$ .

The following closure conditions on constraint languages will be convenient in the further discussion.

**Definition 2.3.** A constraint language is

- closed under renaming iff every constraint has a  $\rho$ -variant for every renaming  $\rho$ ,
- closed under intersection iff for every two constraints  $\phi$  and  $\phi'$  there exists a constraint  $\psi$  s.t.  $\llbracket \phi \rrbracket^{\mathcal{I}} \cap \llbracket \phi' \rrbracket^{\mathcal{I}} = \llbracket \psi \rrbracket^{\mathcal{I}}$  for every interpretation  $\mathcal{I}$ ,
- decidable iff the satisfiability of its constraints is decidable. A constraint  $\phi$  is satisfiable iff there exists at least one interpretation in which  $\phi$  has a solution.

## 2.2.2 Relationally Extended Constraint Languages

To obtain constraint logic programs, a given constraint language  $\mathcal{L}$  has to be extended to a constraint language  $\mathcal{R}(\mathcal{L})$  providing for the necessary relational atoms and propositional connectives.

**Definition 2.4** ( $\mathcal{R}(\mathcal{L})$ ). A constraint language  $\mathcal{R}(\mathcal{L})$  extending a constraint language  $\mathcal{L}$  is defined as follows:

- The signature of  $\mathcal{R}(\mathcal{L})$  is an extension of the signature of  $\mathcal{L}$  with a decidable set  $\mathcal{R}$  of relation symbols and an arity function  $Ar : \mathcal{R} \to \mathbb{N}$ .
- The variables of  $\mathcal{R}(\mathcal{L})$  are the variables of  $\mathcal{L}$ .
- The set of  $\mathcal{R}(\mathcal{L})$ -constraints is the smallest set s.t.
  - 1.  $\phi$  is an  $\mathcal{R}(\mathcal{L})$ -constraint if  $\phi$  is an  $\mathcal{L}$ -constraint,
  - 2.  $r(\vec{x})$  is an  $\mathcal{R}(\mathcal{L})$ -constraint, called an **atom**, if  $r \in \mathcal{R}$  is a relation symbol with arity n and  $\vec{x}$  is an n-tuple of pairwise distinct variables,
  - 3.  $\emptyset$ , F & G,  $F \to G$  are  $\mathcal{R}(\mathcal{L})$ -constraints, if F and G are  $\mathcal{R}(\mathcal{L})$ -constraints,

- 4.  $\phi \& B_1 \& \dots \& B_n \to A \text{ is an } \mathcal{R}(\mathcal{L})\text{-constraint, called a definite clause, if } A, B_1, \dots, B_n \text{ are atoms and } \phi \text{ is an } \mathcal{L}\text{-constraint. We may write a definite clause also as } A \leftarrow \phi \& B_1 \& \dots \& B_n.$
- The variables constrained by an R(L) -constraint are defined as follows: If φ is an L constraint, then V(φ) is defined as in L; V(r(x<sub>1</sub>,...,x<sub>n</sub>)) := {x<sub>1</sub>,...,x<sub>n</sub>}; V(Ø) := Ø; V(F & G) := V(F) ∪ V(G); V(F → G) := V(F) ∪ V(G).
- For each  $\mathcal{L}$ -interpretation  $\mathcal{I}$ , an  $\mathcal{R}(\mathcal{L})$ -interpretation  $\mathcal{A}$  is an extension of an  $\mathcal{L}$ -interpretation  $\mathcal{I}$  with relations  $r^{\mathcal{A}}$  on the domain  $\mathcal{D}$  of  $\mathcal{A}$  with appropriate arity for every  $r \in \mathcal{R}$ , and the domain of  $\mathcal{A}$  is the domain of  $\mathcal{I}$ .
- For each R(L) -interpretation A, for each L -interpretation I, [[·]]<sup>A</sup> is a function mapping every R(L) -constraint to a set of variable assignments s.t.
  - 1.  $\llbracket \phi \rrbracket^{\mathcal{A}} = \llbracket \phi \rrbracket^{\mathcal{I}}$  if  $\phi$  is an  $\mathcal{L}$ -constraint,
  - 2.  $\llbracket r(\vec{x}) \rrbracket^{\mathcal{A}} = \{ \alpha \in \mathsf{ASS} \mid \alpha(\vec{x}) \in r^{\mathcal{A}} \},\$
  - 3.  $\llbracket \emptyset \rrbracket^{\mathcal{A}} = \mathsf{ASS},$
  - $4. \quad \llbracket F \& G \rrbracket^{\mathcal{A}} = \llbracket F \rrbracket^{\mathcal{A}} \cap \llbracket G \rrbracket^{\mathcal{A}},$
  - 5.  $\llbracket F \to G \rrbracket^{\mathcal{A}} = (\mathsf{ASS} \setminus \llbracket F \rrbracket^{\mathcal{A}}) \cup \llbracket G \rrbracket^{\mathcal{A}}.$

Note that we slightly abuse the notation  $\alpha(\vec{x})$  to abbreviate the notation  $(\alpha(x_1), \alpha(x_2), \ldots, \alpha(x_n))$  for a n-tuple of objects assigned to a n-tuple  $\vec{x}$  of variables by a variable assignment  $\alpha$ .

#### 2.2.3 Syntax and Declarative Semantics of Definite Clause Specifications

The concept of a constraint logic program now can be defined as a definite clause specification over a constraint language.

**Definition 2.5 (Definite clause specification).** A definite clause specification  $\mathcal{P}$  over a constraint language  $\mathcal{L}$  is a set of definite clauses from a constraint language  $\mathcal{R}(\mathcal{L})$  extending  $\mathcal{L}$ .

Models of definite clause specifications are determined by the definite clauses constituting these specifications, i.e., a definite clause specification has its definite clauses as its axioms. For reasons of generality, the following two definitions are made with respect to general sets of  $\mathcal{R}(\mathcal{L})$ -constraints.

**Definition 2.6 (Model).** An  $\mathcal{R}(\mathcal{L})$ -interpretation  $\mathcal{A}$  is a model of a set  $\Psi$  of  $\mathcal{R}(\mathcal{L})$ -constraints iff for every  $\alpha \in ASS$ , for every  $\psi \in \Psi$ :  $\alpha \in \llbracket \psi \rrbracket^{\mathcal{A}}$ .

For convenience we furthermore introduce the concept of logical consequence.

**Definition 2.7 (Logical consequence).** An  $\mathcal{R}(\mathcal{L})$ -constraint  $\psi$  is a logical consequence of set  $\Psi$  of  $\mathcal{R}(\mathcal{L})$ -constraints iff, for every  $\mathcal{R}(\mathcal{L})$ -interpretation  $\mathcal{A}$ ,  $\mathcal{A}$  is a model of  $\Psi$  implies that  $\mathcal{A}$  is a model of  $\psi$ .

A goal G is defined as a possibly empty conjunction of  $\mathcal{L}$ -constraints and  $\mathcal{R}(\mathcal{L})$ -atoms.

Given a definite clause specification  $\mathcal{P}$  and a goal G, a  $\mathcal{P}$ -answer of G is defined as a satisfiable  $\mathcal{L}$ -constraint  $\phi$  such that the implication  $\phi \to G$  is a logical consequence of  $\mathcal{P}$ .

In order to show that the semantic properties of conventional logic programming extend to CLP, Höhfeld and Smolka (1988) first define a partial ordering on the set of  $\mathcal{R}(\mathcal{L})$ interpretations.  $\mathcal{R}(\mathcal{L})$ -interpretations extending the same  $\mathcal{L}$ -interpretation  $\mathcal{I}$  are called base equivalent, and  $\mathcal{I}$  is called the base of these  $\mathcal{R}(\mathcal{L})$ -interpretations. A partial ordering on such  $\mathcal{R}(\mathcal{L})$ -interpretations is defined via a partial ordering on the set of the denotations of the relation symbols in these interpretations. We get for all base equivalent  $\mathcal{R}(\mathcal{L})$ -interpretations  $\mathcal{A}, \mathcal{A}'$ :

- $\mathcal{A} \subseteq \mathcal{A}'$  iff for each n-ary relation symbol  $r \in \mathcal{R} : r^{\mathcal{A}} \subseteq r^{\mathcal{A}'}$ ,
- $\mathcal{A} = \bigcup X$  iff for each n-ary relation symbol  $r \in \mathcal{R} : r^{\mathcal{A}} = \bigcup \{ r^{\mathcal{A}'} | \mathcal{A}' \in X \},$
- $\mathcal{A} = \bigcap X$  iff for each n-ary relation symbol  $r \in \mathcal{R} : r^{\mathcal{A}} = \bigcap \{ r^{\mathcal{A}'} | \mathcal{A}' \in X \}.$

This set of base equivalent  $\mathcal{R}(\mathcal{L})$ -interpretations is a complete lattice under the partial order of set inclusion. That is, for every set of base-equivalent  $\mathcal{R}(\mathcal{L})$ -interpretations we have a supremum, given by the union, and an infimum, given by the intersection of the interpretations in the set. The top element is the  $\mathcal{R}(\mathcal{L})$ -interpretation  $\mathcal{A}^{\top}$  such that for each n-ary relation symbol  $r \in \mathcal{R} : r^{\mathcal{A}^{\top}} = \mathcal{D}^{\mathsf{Ar}(r)}$ , and the bottom element is  $\mathcal{A}^{\perp}$  s.t for each n-ary relation symbol  $r \in \mathcal{R} : r^{\mathcal{A}^{\perp}} = \emptyset$ .

Proposition 2.1, due to Höhfeld and Smolka (1988), generalizes the fixpoint- or latticetheoretic semantics of conventional logic programming to CLP. It says that for each  $\mathcal{L}$ interpretation  $\mathcal{I}$ , a definite clause specification  $\mathcal{P}$  in  $\mathcal{R}(\mathcal{L})$  defines unique minimal denotations for the relation symbols of  $\mathcal{R}$ . That is, every  $\mathcal{L}$ -interpretation  $\mathcal{I}$  can be used to construct a minimal model for  $\mathcal{P}$  in  $\mathcal{R}(\mathcal{L})$ . All questions concering the declarative semantics of CLP can then be dealt with in terms of a minimal model semantics. Moreover, a minimal model semantics is crucial for the construction of a sound and complete deduction system for CLP.

**Proposition 2.1 (Höhfeld and Smolka (1988), Theorem 4.4.).** Let  $\mathcal{I}$  be an  $\mathcal{L}$ -interpretation and  $\mathcal{P}$  be a definite clause specification in  $\mathcal{R}(\mathcal{L})$ . Then the equations

$$\begin{aligned} r^{\mathcal{A}_0} &:= \emptyset, \\ r^{\mathcal{A}_{i+1}} &:= \{\alpha(\vec{x}) | \text{ there is a clause } (r(\vec{x}) \leftarrow G) \in \mathcal{P} \text{ and } \alpha \in \llbracket G \rrbracket^{\mathcal{A}_i} \} \end{aligned}$$

- (i) define a chain  $\mathcal{A}_0 \subseteq \mathcal{A}_1 \subseteq \ldots$  of  $\mathcal{R}(\mathcal{L})$ -interpretations extending  $\mathcal{I}$ ,
- (ii) the union  $\mathcal{A} := \bigcup_{i>0} \mathcal{A}_i$  is a model of  $\mathcal{P}$  extending  $\mathcal{I}$ ,
- (iii)  ${\mathcal A} \text{ is the minimal model of } {\mathcal P} \text{ extending } {\mathcal I}$  .

Proposition 2.2 connects the concept of a  $\mathcal{P}$ -answer with the minimal model semantics of  $\mathcal{P}$  (see Höhfeld and Smolka (1988), Proposition 4.5.). This proposition justifies the restriction of the declarative semantics of CLP to a minimal model semantics. We prove this proposition explicitly with reference to the concept of logical consequence.

**Proposition 2.2.** For each definite clause specification  $\mathcal{P}$  in  $\mathcal{R}(\mathcal{L})$ , for each goal G, for each  $\mathcal{L}$ -constraint  $\phi: \phi \to G$  is a logical consequence of  $\mathcal{P}$  iff each minimal model  $\mathcal{A}$  of  $\mathcal{P}$  is a model of  $\phi \to G$ .

*Proof.* If: For each minimal model  $\mathcal{A}$  of  $\mathcal{P}$ :  $\mathcal{A}$  is a model of  $\phi \to G$ 

- $\implies \text{for every model } \mathcal{B} \text{ of } \mathcal{P} \text{ base equivalent to some minimal model } \mathcal{A} \text{ of } \mathcal{P} \text{: } \mathcal{B} \text{ is a model of } \\ \phi \to G, \text{ since } \mathcal{A} \subseteq \mathcal{B} \text{ by Proposition 2.1}$
- $\implies \phi \to G$  is a logical consequence of  $\mathcal{P}$ .
- Only if:  $\phi \to G$  is a logical consequence of  $\mathcal{P}$
- $\implies$  every model of  $\mathcal{P}$  is a model of  $\phi \to G$ , by Definition 2.7
- $\implies \mathcal{A}$  is a model of  $\phi \to G$ .

#### 2.2.4 Operational Semantics of Definite Clause Specifications

The following definitions are made with respect to some implicit  $\mathcal{L}$ ,  $\mathcal{R}$ ,  $\mathcal{P}$ , and  $\mathsf{V}$ , where  $\mathsf{V}$  denotes the finite set of variables in the query and the  $\mathsf{V}$ -solutions of a constraint  $\phi$  in an interpretation  $\mathcal{I}$  are defined as  $\llbracket \phi \rrbracket_{\mathsf{V}}^{\mathcal{I}} := \{ \alpha |_{\mathsf{V}} | \ \alpha \in \llbracket \phi \rrbracket^{\mathcal{I}} \}$  and  $\alpha |_{\mathsf{V}}$  is the restriction of  $\alpha$  to  $\mathsf{V}$ .

Höhfeld and Smolka (1988) define the generalization of the SLD-resolution rule by a binary relation  $\stackrel{r}{\longrightarrow}$ , called **goal reduction**, on the set of goals. The rule selects the leftmost atom in the goal, looks for a variant of a program clause with the selected atom as head, and replaces the selected atom in the goal by the body of the variant clause. Furthermore, the rule ensures that no accidental variable sharing is introduced by the variant.

 $A \& G \xrightarrow{r} F \& G \text{ if } A \leftarrow F \text{ is a variant of a clause in } \mathcal{P}$ s.t.  $(\mathsf{V} \cup \mathsf{V}(G)) \cap \mathsf{V}(F) \subseteq \mathsf{V}(A).$ 

A second rule takes care of **constraint solving** for the  $\mathcal{L}$ -constraints appearing in subsequent goals. The rule takes the conjunction of the  $\mathcal{L}$ -constraints from the reduced goal and the

applied clause and gives, via the black box of a suitable  $\mathcal{L}$ - constraint solver, a satisfiable  $\mathcal{L}$ constraint in solved form if the conjunction of  $\mathcal{L}$ -constraints is satisfiable. If the conjunction of  $\mathcal{L}$ -constraints is not satisfiable, an  $\mathcal{L}$ -constraint – denoting failure is returned. The constraint
solving rule can then be defined as a total function  $\stackrel{c}{\longrightarrow}$  on the set of goals.

 $\phi \And \phi' \And G \xrightarrow{c} \phi'' \And G \text{ if } \llbracket \phi \And \phi' \rrbracket^{\mathcal{I}}_{\mathsf{V} \cup \mathsf{V}(G)} = \llbracket \phi'' \rrbracket^{\mathcal{I}}_{\mathsf{V} \cup \mathsf{V}(G)}$ for all  $\mathcal{L}$ -interpretations  $\mathcal{I}$  and for all  $\mathcal{L}$ -constraints  $\phi, \phi'$  and  $\phi''$ .

Furthermore, a complexity measure that mirrors the construction steps of a minimal model in the complexity of goal reduction is introduced. This measure will be crucial for proving completeness of goal reduction.

• The complexity of a variable assignment  $\alpha$  for an atom A in the minimal model  $\mathcal{A}$  where  $\alpha \in \llbracket A \rrbracket^{\mathcal{A}}$  is defined as

$$comp(\alpha, A, \mathcal{A}) := \min\{i \mid \alpha \in \llbracket A \rrbracket^{\mathcal{A}_i}\};$$

• The complexity of  $\alpha$  for goal G in  $\mathcal{A}$  where  $\alpha \in \llbracket G \rrbracket^{\mathcal{A}}$  is

 $comp(\alpha, G, \mathcal{A}) := \{ comp(\alpha, A, \mathcal{A}) \mid A \text{ is an atom in } G \}$ 

where  $\{\ldots\}$  is a multiset;

• The V-complexity of  $\alpha$  for G in  $\mathcal{A}$  where  $\alpha \in \llbracket G \rrbracket_{\mathsf{V}}^{\mathcal{A}}$  is

 $comp_{\mathsf{V}}(\alpha, G, \mathcal{A}) := \min\{comp(\beta, G, \mathcal{A}) | \beta \in \llbracket G \rrbracket^{\mathcal{A}} \text{ and } \alpha = \beta|_{\mathsf{V}} \}$ 

where  $\beta|_{\mathsf{V}}$  is the restriction of  $\beta$  to the variables in  $\mathsf{V}$ , and the minimum is taken with respect to a total ordering on multisets such that  $M \leq M'$  iff  $\forall x \in M \setminus M', \exists x' \in M' \setminus M$ s.t. x < x'.

Höhfeld and Smolka (1988) prove the following propositions showing that goal reduction is a sound and complete rule for deducing  $\mathcal{P}$ -answers from general definite clause specifications. We prove the main results explicitly in Propositions 2.4 (soundness) and 2.6 (completeness). Note that soundness and completeness can be proven without reference to constraint solving,

**Proposition 2.3 (Höhfeld and Smolka (1988), Proposition 5.1.).** If  $G_1 \xrightarrow{r} G_2$ , then  $\llbracket G_2 \rrbracket^{\mathcal{A}} \subseteq \llbracket G_1 \rrbracket^{\mathcal{A}}$  for every model  $\mathcal{A}$  of  $\mathcal{P}$ .

**Proposition 2.4.** If  $G \xrightarrow{r} \phi$ , then  $\phi \to G$  is a logical consequence of  $\mathcal{P}$ .

*Proof.*  $G \xrightarrow{r} \phi^* \phi$ 

 $\implies \llbracket \phi \rrbracket^{\mathcal{A}} \subseteq \llbracket G \rrbracket^{\mathcal{A}}$  for every model  $\mathcal{A}$  of  $\mathcal{P}$ , by Proposition 2.3 and transitivity of  $\subseteq$ 

- $\implies \text{for every model } \mathcal{A} \text{ of } \mathcal{P}: \llbracket \phi \to G \rrbracket^{\mathcal{A}} = \mathsf{ASS}, \text{ since for every model } \mathcal{A} \text{ of } \mathcal{P}: \llbracket \phi \rrbracket^{\mathcal{A}} \subseteq \llbracket G \rrbracket^{\mathcal{A}}$
- $\implies$  for every model  $\mathcal{A}$  of  $\mathcal{P}$ :  $\mathcal{A}$  is a model of  $\phi \to G$ , by Definition 2.6
- $\implies \phi \to G$  is a logical consequence of  $\mathcal{P}$ .

**Proposition 2.5 (Höhfeld and Smolka (1988), Theorem 5.2.).** Let  $\mathcal{L}$  be closed under renaming,  $\mathcal{A}$  be a minimal model of  $\mathcal{P}$ ,  $G_1$  be a goal, A be an atom in  $G_1$ , and  $\alpha \in \llbracket G_1 \rrbracket_V^{\mathcal{A}}$ . Then there exists a clause C in  $\mathcal{P}$  and a goal  $G_2$  s.t.  $G_1 \xrightarrow{r} G_2$  using a variant of C on A is possible,  $\alpha \in \llbracket G_2 \rrbracket_V^{\mathcal{A}}$  and  $comp_V(\alpha, G_2, \mathcal{A}) < comp_V(\alpha, G_1, \mathcal{A})$ .

**Proposition 2.6 (Höhfeld and Smolka (1988), Corollary 5.3.).** Let  $\mathcal{L}$  be closed under renaming,  $\mathcal{A}$  be a minimal model of  $\mathcal{P}$ , G be a goal and  $\alpha \in \llbracket G \rrbracket_{\mathsf{V}}^{\mathcal{A}}$ . Then there exists a  $\mathcal{P}$ -answer  $\phi$  of G s.t.  $G \xrightarrow{r} \phi$  and  $\alpha \in \llbracket \phi \rrbracket_{\mathsf{V}}^{\mathcal{A}}$ .

- *Proof.* The result is proven by induction on  $comp_{V}(\alpha, G, \mathcal{A})$ .
- Base: Goals with mulitset complexity  $\emptyset$  have to be a satisfiable  $\mathcal{L}$ -constraint  $\phi$ . Then  $\phi \xrightarrow{r}{\longrightarrow} {}^{0}\phi$  and  $\phi$  is a  $\mathcal{P}$ -answer of itself.
- Hypothesis: Suppose the result holds for goals with multiset complexity less than some multiset N.

Step:  $comp_{\mathsf{V}}(\alpha', G_1, \mathcal{A}) = N$  and  $\alpha' \in \llbracket G_1 \rrbracket_{\mathsf{V}}^{\mathcal{A}}$ 

- $\implies$  there exists a clause C of  $\mathcal{P}$  and a goal  $G_2$  s.t.  $G_1 \xrightarrow{r} G_2$  and  $\alpha' \in \llbracket G_2 \rrbracket_{\mathsf{V}}^{\mathcal{A}}$  and  $comp_{\mathsf{V}}(\alpha', G_2, \mathcal{A}) < comp_{\mathsf{V}}(\alpha', G_1, \mathcal{A})$ , by Proposition 2.5
- $\implies$  there exists a  $\mathcal{P}$ -answer  $\phi$  of  $G_2$  s.t.  $G_2 \xrightarrow{r} \phi$  and  $\alpha' \in \llbracket \phi \rrbracket_V^{\mathcal{A}}$ , by the hypothesis
- $\implies$  there exists a  $\mathcal{P}$ -answer  $\phi$  of  $G_1$  s.t.  $G_1 \xrightarrow{r} \phi^* \phi$  and  $\alpha' \in \llbracket \phi \rrbracket_V^{\mathcal{A}}$ , and by Proposition 2.4,  $\phi \to G_1$  is a logical consequence of  $\mathcal{P}$ .

The result follows by arithmetic induction.

In all following examples, we will use a standard Prolog resolution procedure for the CLP scheme of Höhfeld and Smolka (1988), i.e., we combine the left-right selection rule defined in goal reduction with a depth-first search rule. Furthermore, after each goal reduction step, constraint solving is applied, and another clause is tried immediately if constraint solving fails. Moreover, it will be convenient in the following discussion to view the search space determined by the derivation rules  $\xrightarrow{r}$  and  $\xrightarrow{c}$  as a search of a tree. A derivation tree is defined as follows.

**Definition 2.8 (Derivation tree).** A derivation tree determined by a query  $G_1$  and a definite clause specification  $\mathcal{P}$  has to satisfy the following conditions:

1. Each node is either a relation node or a constraint node.

- The successors of every relation node are all constraint nodes s.t. for every -<sup>r</sup>→ -resolvent G' obtainable by a clause C from goal G in a relation node, there is a successor constraint node labeled by C and G'.
- 3. The successors of every constraint node are all relation nodes s.t. for the unique  $\xrightarrow{c}$  resolvent  $G\&\phi''$  obtainable from goal  $G\&\phi\&\phi'$  in a constraint node, there is a successor relation node labeled by  $G\&\phi''$ .
- 4. The root node is a relation node labeled by  $G_1$ .
- 5. A success node is a terminal relation node labeled by a satisfiable  $\mathcal{L}$  -constraint.

Successful derivations correspond to subtrees of derivation trees which are labeled by terminal success nodes. Such trees can be defined as proof trees as follows.

**Definition 2.9 (Proof tree).** A proof tree for a query  $G_1$  from  $\mathcal{P}$  is a subtree of a derivation tree determined by  $G_1$  and  $\mathcal{P}$  and is defined as follows:

- 1. A relation node of the proof tree is a relation node of the supertree and takes one of the successors of the relation node of the supertree as its successor node.
- 2. A constraint node of the proof tree is a constraint node of the supertree and takes the unique successor of the constraint node of the supertree as its successor node.
- 3. The root node of the proof tree is the root node of the supertree.
- The terminal node of the proof tree is a success node of the supertree, labeled by a satisfiable L -constraint, called answer constraint.

Let us illustrate the basic concepts of CLP with an example. A simple program consisting of clauses 1 to 3 is depicted in Fig. 2.1.

$$1 q(X) \leftarrow p(X).$$
  

$$2 p(X) \leftarrow X = a.$$
  

$$3 p(X) \leftarrow X = b.$$

Figure 2.1: Constraint logic program

The  $\mathcal{L}$ -constraints are considered to come from a language of hierarchical types, where the ordering on types is defined via the operation of set inclusion on their denotations. In our example, we have  $\llbracket a \rrbracket^{\mathcal{I}} \subseteq \llbracket e \rrbracket^{\mathcal{I}}$ ,  $\llbracket b \rrbracket^{\mathcal{I}} \subseteq \llbracket e \rrbracket^{\mathcal{I}}$  and  $\llbracket a \rrbracket^{\mathcal{I}} \cap \llbracket b \rrbracket^{\mathcal{I}} = \emptyset$ . This hierarchy is depicted graphically in Fig. 2.2.



Figure 2.2: Type hierarchy

$$\begin{aligned} \mathbf{p}^{\mathcal{A}_{0}} &= \emptyset, \mathbf{q}^{\mathcal{A}_{0}} = \emptyset, \\ \mathbf{p}^{\mathcal{A}_{1}} &= \{ \llbracket a \rrbracket^{\mathcal{I}}, \llbracket b \rrbracket^{\mathcal{I}} \}, \mathbf{q}^{\mathcal{A}_{1}} = \emptyset, \\ \mathbf{p}^{\mathcal{A}_{2}} &= \{ \llbracket a \rrbracket^{\mathcal{I}}, \llbracket b \rrbracket^{\mathcal{I}} \}, \mathbf{q}^{\mathcal{A}_{2}} = \{ \llbracket a \rrbracket^{\mathcal{I}}, \llbracket b \rrbracket^{\mathcal{I}} \}, \\ \vdots \\ \mathbf{p}^{\mathcal{A}} &= \{ \llbracket a \rrbracket^{\mathcal{I}}, \llbracket b \rrbracket^{\mathcal{I}} \}, \mathbf{q}^{\mathcal{A}} = \{ \llbracket a \rrbracket^{\mathcal{I}}, \llbracket b \rrbracket^{\mathcal{I}} \}, \text{ where } \mathcal{A} = \bigcup_{i \geq 0} \mathcal{A}_{i}. \end{aligned}$$

Figure 2.3: Minimal model construction for constraint logic program

$$q(X) \& X = e \\ | r \\ 1, p(X) \& X = e \\ | c \\ p(X) \& X = e \\ r \\ r \\ 2, X = e \& X = a \\ | c \\ x = a \\ X = b \\$$

Figure 2.4: Derivation tree for constraint logic program

The construction of a minimal model for the program of Fig. 2.1 is shown in Fig. 2.3. The unique minimal denotations of the relation symbols p and q are obtained in step 1 and 2 of the minimal model construction respectively.

A derivation tree for the query q(X) & X = e from the program of Fig. 2.1 is given in Fig. 2.4. We depict only the success branches of the derivation tree, yielding two distinct proof trees for the query, with answer constraints X = a and X = b respectively.

Soundness of the CLP scheme implies that corresponding to the derivation of X = a and X = b, we know that the implications  $X = a \rightarrow q(X)$  & X = e and  $X = b \rightarrow q(X)$  & X = e are logical consequences of the program of Fig. 2.1. This is easily verified from the minimal model given in Fig. 2.3. Furthermore, completeness of the CLP scheme is easily verified from the fact that for solutions  $\alpha \in [\![q(X) \& X = e]\!]_V^A$  and  $\alpha' \in [\![q(X) \& X = e]\!]_V^A$ , we can derive  $\mathcal{P}$ -answers with  $\alpha \in [\![X = a]\!]_V^A$  and  $\alpha' \in [\![X = b]\!]_V^A$ .

## 2.3 Constraint Logic Grammars

In this section, we will explicate the concept of constraint logic grammars. To this end, we will restrict our attention to feature-based CLGs and discuss in particular the main properties of an HPSG instance of such grammars.

We will show how a feature-based constraint language can be obtained from a featurebased logical description language, and how such a constraint language can be embedded into the CLP scheme of Höhfeld and Smolka (1988), yielding a feature-based CLG. The language to be discussed is that of Götz (1995), Götz (to appear), which is close to that of King (1989), King (1994) (modulo the usage of variables) and Smolka (1988), Smolka (1992) (modulo appropriateness conditions). This language provides a description language  $\mathcal{FD}$  specifying the logical foundations of HPSG grammars and is extendable to a constraint language  $\mathcal{FL}$ , in the sense of Höhfeld and Smolka (1988). The expressive power of the language is smaller than or equal to the expressive power of first-order predicate logic with equality.

#### 2.3.1 A Feature-Based Constraint Language

The language is based on a notion of signature, i.e., the non-logical elements of the alphabet, declaring the structures the linguist is interested in. A signature specifies a set of feature symbols, a lattice of sort symbols and appropriateness conditions restricting the functional properties of the feature symbols. All subsequent work should be understood with respect to an implicit signature  $\Sigma$ .

**Definition 2.10 (Signature).** A signature is a quadruple  $\langle \mathcal{T}, \preceq, \mathcal{F}, approp \rangle$  s.t.

•  $\langle \mathcal{T}, \preceq \rangle$  is a finite join-semilattice of types,

- $S = \{t \in \mathcal{T} | if t' \leq t then t' = t\}$  is a finite set of minimal types,
- $\mathcal{F}$  is a finite set of feature symbols,
- approp :  $S \times F \rightarrow T$  is a partial function from pairs of minimal types and features to types.

The well-formed formulae of the feature-based description language  $\mathcal{FD}$ , called feature descriptions, are built from the symbols in the signature, a countably infinite set of variables VAR, the symbol : assigning features to their values, and the standard boolean connectives. Expressions of this kind can be seen as the formal equivalent of the AVM notation used in Pollard and Sag (1994). The set **Desc** of feature descriptions is defined as follows.

**Definition 2.11 (Feature descriptions).** The set Desc of feature descriptions is the smallest set s.t.

• X is a description	$if X \in VAR,$
• t is a description	if $t \in \mathcal{T}$ ,
• f:D is a description	$\textit{if}f\in\mathcal{F},\;D\inDesc,$

•  $D_1 \wedge D_2$ ,  $D_1 \vee D_2$ ,  $\neg D_1$ ,  $D_1 \rightarrow D_2$  are descriptions if  $D_1 \in \mathsf{Desc}$ ,  $D_2 \in \mathsf{Desc}$ .

An interpretation of a signature is based on an arbitrary domain of objects, and assigns to every object exactly one minimal type, and to every feature symbol a partial function on the domain. The domains and ranges of these functions are determined by the *approp* function. This function specifies that for each object u of a minimal type s, there is a connected object F(f)(u) defined iff approp(s, f) is defined, and the type S(F(f)(u)) of this connected object has to be appropriate.

**Definition 2.12 (Interpretation).** An interpretation is a quadruple  $\mathcal{I} = \langle U, S, F \rangle$  s.t.

- U is a set of objects, the domain of  $\mathcal{I}$ ,
- S:  $U \rightarrow S$  is a total function from the domain to the set of minimal types,
- $F: \mathcal{F} \to U^U$  is a is a total feature interpretation function s.t.
  - 1. for each  $u \in U$ , for each  $f \in \mathcal{F}$ , if approp(S(u), f) is defined and approp(S(u), f) = t, then F(f)(u) is defined and  $S(F(f)(u)) \leq t$ ,
  - 2. for each  $u \in U$ , for each  $f \in \mathcal{F}$ , if F(f)(u) is defined, then approp(S(u), f) is defined and  $S(F(f)(u)) \preceq approp(S(u), f)$ .
The denotation of feature descriptions with respect to an interpretation  $\mathcal{I}$  and a variable assignment  $\alpha$  is defined to be a subset of the domain for every feature description. By abstracting away from the variable assignment, we arrive at a concept of abstract denotation comprising the denotation of a feature description under every possible variable assignment.

**Definition 2.13 (Variable assignment).** A variable assignment  $\alpha$  : VAR  $\rightarrow$  U is a total function from the set of variables to the domain. Write ASS for the set of variable assignments.

Definition 2.14 (Feature description denotation).

- $\llbracket X \rrbracket_{\alpha}^{\mathcal{I}} = \{ \alpha(X) \}$  if  $X \in \mathsf{VAR}$ ,
- $\llbracket t \rrbracket_{\alpha}^{\mathcal{I}} = \{ u \in \mathsf{U} | \mathsf{S}(u) \preceq t \}$  if  $t \in \mathcal{T}$ ,
- $\llbracket f: D \rrbracket_{\alpha}^{\mathcal{I}} = \{ u \in \mathsf{U} | \mathsf{F}(f)(u) \text{ is defined, } \mathsf{F}(f)(u) \in \llbracket D \rrbracket_{\alpha}^{\mathcal{I}} \} \text{ if } f \in \mathcal{F}, D \in \mathsf{Desc},$
- $\llbracket D_1 \wedge D_2 \rrbracket_{\alpha}^{\mathcal{I}} = \llbracket D_1 \rrbracket_{\alpha}^{\mathcal{I}} \cap \llbracket D_2 \rrbracket_{\alpha}^{\mathcal{I}}$  if  $D_1, D_2 \in \mathsf{Desc}$ ,
- $\llbracket D_1 \vee D_2 \rrbracket_{\alpha}^{\mathcal{I}} = \llbracket D_1 \rrbracket_{\alpha}^{\mathcal{I}} \cup \llbracket D_2 \rrbracket_{\alpha}^{\mathcal{I}}$  if  $D_1, D_2 \in \mathsf{Desc}$ ,
- $\llbracket \neg D_1 \rrbracket_{\alpha}^{\mathcal{I}} = \mathsf{U} \setminus \llbracket D_1 \rrbracket_{\alpha}^{\mathcal{I}}$  if  $D_1 \in \mathsf{Desc}$ ,
- $\llbracket D_1 \to D_2 \rrbracket_{\alpha}^{\mathcal{I}} = (\mathsf{U} \setminus \llbracket D_1 \rrbracket_{\alpha}^{\mathcal{I}}) \cup \llbracket D_2 \rrbracket_{\alpha}^{\mathcal{I}}$  if  $D_1, D_2 \in \mathsf{Desc.}$

Definition 2.15 (Abstract denotation).

$$\llbracket D \rrbracket^{\mathcal{I}} = \bigcup_{\alpha \in \mathsf{ASS}} \llbracket D \rrbracket^{\mathcal{I}}_{\alpha} \qquad \text{ if } D \in \mathsf{Desc.}$$

To obtain a feature-based constraint language  $\mathcal{FL}$  fulfilling the closure requirements on constraint languages stated by Höhfeld and Smolka (1988), first we simply have to attach every feature description D in  $\mathcal{FD}$  with a new variable not occuring in the set V(D) of variables in D. This avoids accidental variable sharing and guarantees renaming closure of  $\mathcal{FL}$ . Furthermore, an explicit definition of conjunction of feature constraints ensures intersection closure of  $\mathcal{FL}$ .

#### Definition 2.16 (Feature constraints).

- X = D is a constraint if  $X \in VAR$ ,  $X \notin V(D)$ ,  $D \in Desc$ ,
- $\phi \& \phi'$  is a constraint if  $\phi$ ,  $\phi'$  are constraints.

The denotation of a constraint is defined by a function mapping every constraint to a set of variable assignments, called solutions. The solutions of a constraint X = D are the variable assignments in ASS which constrain the value of the variable X to the objects in the denotation of D. The denotation of a conjunction of constraints is the intersection of the respective denotations.

#### Definition 2.17 (Feature constraint solutions).

•  $\llbracket X = D \rrbracket^{\mathcal{I}} = \{ \alpha \in \mathsf{ASS} | \ \alpha(X) \in \llbracket D \rrbracket^{\mathcal{I}}_{\alpha} \}$  if  $X \in \mathsf{VAR}, \ X \notin \mathsf{V}(D), \ D \in \mathsf{Desc},$ •  $\llbracket \phi \And \phi' \rrbracket^{\mathcal{I}} = \llbracket \phi \rrbracket^{\mathcal{I}} \cap \llbracket \phi' \rrbracket^{\mathcal{I}}$  if  $\phi, \ \phi' \ are \ constraints.$ 

Next we have to consider the problem of deciding satisfiability of feature descriptions and feature constraints.

**Definition 2.18 (Satisfiability of feature descriptions).** A feature description D is satisfiable iff there is an interpretation  $\mathcal{I}$  s.t.  $[D]^{\mathcal{I}} \neq \emptyset$ .

This problem has been shown to be decidable for feature-based description languages closely related to the above reported one. For the description language reported above, a decision algorithm is given by Götz (to appear), for the variable-free notational variant of King (1994) by Kepser (1994), for a less expressive version of the language not employing appropriateness conditions by Smolka (1988), Smolka (1992), or for an even less expressive version employing conjunction as only boolean operator by Aït-Kaci, Podelski, and Goldstein (1993).

Most of these approaches adapt for satisfiability checking a constraint solving method similar to that of Smolka (1988), Smolka (1992). This method is a three-step transformation process from feature descriptions to a solved form of feature constraints displaying (un)satisfiability. Following Götz (to appear), constraint solving for the feature-based constraint language reported above can be illustrated as follows: Firstly, every feature description is transformed to disjunctive normal form; secondly, every feature description in disjunctive normal form is transformed into a (disjunctively interpreted) set of (conjunctively interpreted) sets of feature constraints of the simple form X = Y,  $X = \neg Y$ , X = t or X = f : Y; thirdly, every such set of sets of simple feature constraints is transformed into a set of sets of feature constraints in solved form.

For reasons of readability, we will consider the constraint solver for the feature-based constraint language  $\mathcal{FL}$  in the following as a black box. The interested reader is referred for details and proofs to Götz (to appear). In all subsequent examples, we will depict only the result of constraint solving, re-translated from simple feature constraints in solved normal form to feature constraints in a more readable form according to Definition 2.16.

The notion of satisfiability defined for feature constraints is as follows.

**Definition 2.19 (Satisfiability of feature constraints).** A feature constraint  $\phi$  is satisfiable iff there exists an interpretation  $\mathcal{I}$  s.t.  $\llbracket \phi \rrbracket^{\mathcal{I}} \neq \emptyset$ .

Since every feature constraint is satisfiable whenever the embedded feature description is satisfiable, and since satisfiability of feature descriptions is decidable, we get immediately the desired decidability result for the feature-based constraint language  $\mathcal{FL}$ . To sum up, since

 $\mathcal{FL}$  is closed under renaming and intersection, and due to the decidability algorithm for  $\mathcal{FL}$  constraint solving of Götz (to appear), we can state the following proposition.

**Proposition 2.7.**  $\mathcal{FL}$  is a decidable constraint language closed under renaming and intersection.

#### 2.3.2 Feature-Based Constraint Logic Grammars

Feature-based grammars can be built in a pure declarative way simply as sets of axiomatic interpreted feature descriptions from the feature description language  $\mathcal{FD}$ .

**Definition 2.20 (Grammar).** A feature-based grammar  $\mathcal{G}$  is a finite set of feature descriptions s.t.  $\mathcal{G} \subseteq \mathsf{Desc.}$ 

The feature descriptions comprising a grammar constrain the admissible models of the grammar in that in every model of a grammar every feature description must be true of every object.

**Definition 2.21 (Model).** A model of a feature constraint grammar  $\mathcal{G}$  is an interpretation  $\mathcal{I} = \langle \mathsf{U}, \mathsf{S}, \mathsf{F} \rangle$  s.t. for every  $u \in \mathsf{U}$ , for every  $D \in \mathcal{G}$ :  $u \in \llbracket D \rrbracket^{\mathcal{I}}$ .

The central problem of prediction can then be defined model-theoretically as a relation between grammars and feature descriptions encoding the questioned input.

**Definition 2.22 (Prediction).** A feature description D is predicted by a grammar  $\mathcal{G}$  iff there is a model  $\mathcal{I}$  of  $\mathcal{G}$  s.t.  $[\![D]\!]^{\mathcal{I}} \neq \emptyset$ .

In contrast to this definition, the linguistic problem of grammaticality is sometimes considered as a relation between grammars and objects. As we will see below, the syntactic coding of Def. 2.22 enables a connection of the model-theoretic concept of prediction with the implementational parsing/generation problem. The problem of prediction has shown to be undecidable for various feature-based description languages (see Aït-Kaci, Podelski, and Goldstein (1993), Smolka (1992), Götz (to appear)). As shown by Götz (to appear), decidable fragments of such languages are obtainable, e.g., in the form of grammars fulfilling the finite model property.

**Definition 2.23.** A grammar  $\mathcal{G}$  has the finite model property iff for all descriptions D,

 $\mathcal{G}$  predicts D iff  $\mathcal{G}$  has a finite model  $\mathcal{I}$  s.t  $\llbracket D \rrbracket^{\mathcal{I}} \neq \emptyset$ .

Note that even if for grammars having the finite model property the prediction problem is decidable, it is undecidable if a grammar has the finite model property of not. Thus it has to be kept in mind that decidability of the prediction problem for linguistically interesting CLGs is based on an assumption of finiteness of linguistic structures.

To obtain feature-based CLGs from feature-based grammars, the feature descriptions from  $\mathcal{FD}$  have to be extended to feature constraints from  $\mathcal{FL}$ , which then can be embedded as  $\mathcal{FL}$ constraints into a suitable definite clause specification in  $\mathcal{R}(\mathcal{FL})$ . An example for such an
embedding of a feature-based grammar into the CLP scheme of Höhfeld and Smolka (1988)
is given below. The resulting feature-based CLG can be seen as a notational variant of a
CUF-grammar (Dörre and Eisele 1991; Dörre and Dorna 1993). Alternatively, when replacing
the predicates of this feature-based CLG by a single predicate gram in all clauses, we arrive
at a program which would result from a direct application of the compilation algorithm of
Götz (1995), Götz and Meurers (1995) to a feature-based grammar<sup>3</sup>. This compilation scheme
connects the model-theoretic concept of prediction with the logic programming concept of  $\mathcal{P}$ -answer directly. This is done by an automatic generation of a  $\mathcal{R}(\mathcal{FL})$ -program  $\mathcal{P}$  for every  $\mathcal{FD}$ -grammar  $\mathcal{G}$ , where the program defines an unary relation gram encoding prediction. This
encoding is said to be correct under the following conditions.

**Definition 2.24.** Let  $\mathcal{P}$  be a definite clause specification in  $\mathcal{FL}$  defining the relation gram, and let  $\mathcal{G}$  a grammar from  $\mathcal{FD}$ . Then  $\mathcal{P}$  is a correct translation of  $\mathcal{G}$  iff

 $\mathcal{G}$  predicts feature description D iff the goal gram(X) & X = D has a  $\mathcal{P}$ -answer.

The compilation scheme presented by Götz (1995) is sound and for a large class of grammars complete. A sufficient condition to receive correct translations in the sense of Def. 2.24 is again the finite model property. Thus under the assumption that linguistic structures are finite, CLP can be seen as a useful parsing scheme for linguistically interesting feature-based CLGs.

Let us illustrate these concepts with an example. Suppose a simple grammar licensing, among others, analyses such as

[Peter believes 
$$[Clinton_N \ talks_V]_S$$
]<sub>S</sub>

or

```
[Peter believes [Clinton<sub>N</sub> talks<sub>N</sub>]<sub>NP</sub>]<sub>S</sub>.
```

We will define now a feature-based grammar presenting a  $\mathcal{FD}$ -encoding of the part of this grammar which is relevant for the structural ambiguity. It is a modified and extended version of an example from Carpenter (1992).

<sup>&</sup>lt;sup>3</sup>Based on a differentiation of types in distinct sets according to whether and how they appear as antecedents of grammar constraints, this compilaton procedure introduces a set of clauses defining the single predicate gram for each such set of types. Actually, for the example given below, this compilation scheme would also produce a clause gram $(X) \leftarrow X = t$  for each minimal type t of the grammar signature which is not the antecedent of a grammar description. For ease of readability, we will omit clauses introduced for (minimal or non-minimal) non-antecedent types in our example.

The signature comes with a type hierarchy with top element  $\top$ , feature symbols, and appropriateness conditions, and is depicted in the graph in Fig. 2.5. Feature symbols are depicted in SMALL CAPS font, type symbols in *lower case italics*, and appropriateness conditions are expressed in a matrix notation, reading, e.g., approp(phrase, DTR1) = sign.



Figure 2.5: Signature for feature-based grammar

The relevant  $\mathcal{FD}$ -descriptions are given in Fig. 2.6. The first implication encodes the rules  $S \to N V$  and  $NP \to N N$ . Context-sensitivity is introduced by the agreement requirement on the first rule. The second implication encodes the rules  $N \to Clinton, V \to talks$ , and  $N \to talks$ .

$$\begin{array}{ll} phrase \rightarrow & (\mathrm{CAT}:s \,\wedge\, \mathrm{DTR1:CAT}:n \,\wedge\, \mathrm{DTR2:CAT}:v \,\wedge\, \mathrm{DTR1:AGR}:Y) \\ & \wedge\, \mathrm{DTR2:AGR}:Y) \\ & \vee\, (\mathrm{CAT}:np \,\wedge\, \mathrm{DTR1:CAT}:n \,\wedge\, \mathrm{DTR2:CAT}:n) \\ word \rightarrow & (\mathrm{CAT}:n \,\wedge\, \mathrm{PHON}:Clinton \,\wedge\, \mathrm{AGR}:sg) \\ & \vee\, (\mathrm{CAT}:v \,\wedge\, \mathrm{PHON}:talks \,\wedge\, \mathrm{AGR}:sg) \\ & \vee\, (\mathrm{CAT}:n \,\wedge\, \mathrm{PHON}:talks \,\wedge\, \mathrm{AGR}:sg) \\ & \vee\, (\mathrm{CAT}:n \,\wedge\, \mathrm{PHON}:talks \,\wedge\, \mathrm{AGR}:pl) \end{array}$$

Figure 2.6: Feature-based grammar

The CLG obtained from a simplified compilation of the grammar in Fig. 2.6 to a definite clause specification in  $\mathcal{FL}$  is given in Fig. 2.7. The embedded  $\mathcal{FL}$ -constraints are depicted graphically in the same way as  $\mathcal{FD}$ -descriptions.  $\mathcal{R}(\mathcal{FL})$ -atoms are depicted in typewriter font.

Given this program and a goal

 $X = (sign \land DTR1: PHON : Clinton \land DTR2: PHON : talks) \& sign(X)$ 

encoding the phrase *Clinton talks*, we can infer two answers

- 1 phrase(X)  $\leftarrow X = (phrase \land CAT : s \land DTR1:CAT : n \land DTR2:CAT : v \land DTR1:AGR : Y \land DTR2:AGR : Y \land DTR1 : Z_1 \land DTR2 : Z_2) \& \operatorname{sign}(Z_1) \& \operatorname{sign}(Z_2).$
- 2 phrase(X)  $\leftarrow$  X = (phrase  $\wedge$  CAT :  $np \wedge$  DTR1:CAT :  $n \wedge$  DTR2:CAT :  $n \wedge$  DTR1 :  $Z_1 \wedge$  DTR2 :  $Z_2$ ) & sign( $Z_1$ ) & sign( $Z_2$ ).
- $3 \operatorname{word}(X) \leftarrow X = (word \wedge \operatorname{CAT} : n \wedge \operatorname{PHON} : Clinton \wedge \operatorname{AGR} : sg).$
- 4 word(X)  $\leftarrow$  X = (word  $\land$  CAT :  $v \land$  PHON : talks  $\land$  AGR : sg).
- 5 word(X)  $\leftarrow$  X = (word  $\land$  CAT :  $n \land$  PHON : talks  $\land$  AGR : pl).

```
\texttt{6 sign}(X) \leftarrow \texttt{phrase}(X).
```

```
7 \, \operatorname{sign}(X) \gets \operatorname{word}(X).
```

Figure 2.7: Feature-based constraint logic grammar

$$\begin{split} X &= (phrase \ \land \ \mathrm{CAT} : s \ \land \ \mathrm{DTR1} : word \ \land \ \mathrm{DTR1} : \mathrm{CAT} : n \\ \land \ \mathrm{DTR1} : \ \mathrm{PHON} : Clinton \ \land \ \mathrm{DTR1} : \ \mathrm{AGR} : Y \ \land \ \mathrm{DTR1} : \ \mathrm{AGR} : sg \\ \land \ \mathrm{DTR2} : word \ \land \ \mathrm{DTR2} : \ \mathrm{CAT} : v \ \land \ \mathrm{DTR2} : \ \mathrm{PHON} : talks \\ \land \ \mathrm{DTR1} : \ \mathrm{AGR} : Y \ \land \ \mathrm{DTR1} : \ \mathrm{AGR} : sg) \end{split}$$

and

 $X = (phrase \land CAT : np \land DTR1 : word \land DTR1 : CAT : n \land DTR1 : PHON : Clinton \land DTR1 : AGR : sg \land DTR2 : word \land DTR2 : CAT : n \land DTR2 : PHON : talks \land DTR2 : AGR : pl)$ 

encoding the parses  $[Clinton_N \ talks_V]_S$  and  $[Clinton_N \ talks_N]_{NP}$  respectively. The parses are depicted in Figs. 2.8 and 2.9. Note that goal reduction and constraint solving are applied in one step. Furthermore, only success branches are depicted and the the constraint solver is viewed as a black box.

# 2.4 Summary

In this chapter we discussed the basic formal concepts of the CLP scheme of Höhfeld and Smolka (1988). These concepts provide a formal specification of the notions of constraint language and of a constraint logic program embedding a constraint language. For convenience, we gave some missing proofs and introduced the notions of logical consequence, derivation tree and proof tree into the CLP scheme. These concepts will be useful in latter chapters.

 $X = (sign \land DTR1: PHON : Clinton \land DTR2: PHON : talks)$  $\& \operatorname{sign}(X)$ r, c6,  $X = (sign \land DTR1: PHON : Clinton \land DTR2: PHON : talks)$ & phrase(X) r, c1,  $X = (phrase \land CAT : s \land DTR1 : word \land DTR1 : CAT : n$  $\wedge$  DTR1: PHON : Clinton  $\wedge$  DTR1: AGR : Y  $\wedge$  DTR2 : word  $\wedge$  DTR2: CAT : v  $\wedge$  DTR2: PHON : talks  $\wedge$  DTR2: AGR :  $Y \wedge$  DTR1 :  $Z_1 \wedge$  DTR2 :  $Z_2$ ) & sign $(Z_1)$  & sign $(Z_2)$ r, c7,  $X = (phrase \land CAT : s \land DTR1 : word \land DTR1 : CAT : n$  $\wedge$  DTR1: PHON : Clinton  $\wedge$  DTR1: AGR : Y  $\wedge$  DTR2 : word  $\wedge$  DTR2: CAT : v  $\wedge$  DTR2: PHON : talks  $\wedge$  DTR2: AGR :  $Y \wedge$  DTR1 :  $Z_1 \wedge$  DTR2 :  $Z_2$ )  $\& \operatorname{word}(Z_1) \& \operatorname{sign}(Z_2)$ r, c3,  $X = (phrase \land CAT : s \land DTR1 : word \land DTR1 : CAT : n$  $\wedge$  DTR1: PHON : Clinton  $\wedge$  DTR1: AGR :  $Y \wedge$  DTR1: AGR :  $sq \wedge$  DTR2 : word  $\wedge$  DTR2: CAT :  $v \wedge$  DTR2: PHON : talks  $\wedge$  DTR2: AGR :  $Y \wedge$  DTR2: AGR : sq  $\wedge \operatorname{DTR2}: Z_2) \,\&\, \operatorname{\mathtt{sign}}(Z_2)$ r, c7,  $X = (phrase \land CAT : s \land DTR1 : word \land DTR1 : CAT : n$  $\wedge$  DTR1: PHON : Clinton  $\wedge$  DTR1: AGR :  $Y \wedge$  DTR1: AGR :  $sg \wedge$  DTR2 : word  $\wedge$  DTR2: CAT :  $v \wedge$  DTR2: PHON : talks  $\wedge$  DTR2: AGR :  $Y \wedge$  DTR2: AGR : sq  $\wedge \operatorname{DTR2}: Z_2)$  & word $(Z_2)$ r, c4,  $X = (phrase \land CAT : s \land DTR1 : word \land DTR1 : CAT : n$  $\wedge$  DTR1: PHON : Clinton  $\wedge$  DTR1: AGR :  $Y \wedge$  DTR1: AGR :  $sq \wedge$  DTR2 : word  $\wedge$  DTR2: CAT :  $v \wedge$  DTR2: PHON : talks  $\wedge$  DTR2: AGR :  $Y \wedge$  DTR2: AGR : sq)

Figure 2.8: A derivation of  $[Clinton_N \ talks_V]_S$ 

 $X = (sign \land DTR1: PHON : Clinton \land DTR2: PHON : talks)$  $\& \operatorname{sign}(X)$ r, c6,  $X = (sign \land DTR1: PHON : Clinton \land DTR2: PHON : talks)$ & phrase(X) r, c2,  $X = (phrase \land CAT : np \land DTR1 : word \land DTR1 : CAT : n$  $\land$  DTR1: PHON : Clinton  $\land$  DTR2 : word  $\land$  DTR2: CAT : n  $\wedge$  DTR2: PHON : talks  $\wedge$  DTR1 :  $Z_1 \wedge$  DTR2 :  $Z_2$ )  $\& \operatorname{sign}(Z_1) \& \operatorname{sign}(Z_2)$ r, c7,  $X = (phrase \land CAT : np \land DTR1 : word \land DTR1 : CAT : n$  $\land$  DTR1: PHON : Clinton  $\land$  DTR2 : word  $\land$  DTR2: CAT : n  $\wedge$  DTR2: PHON : talks  $\wedge$  DTR1 :  $Z_1 \wedge$  DTR2 :  $Z_2$ )  $\& \operatorname{word}(Z_1) \& \operatorname{sign}(Z_2)$ r, c3,  $X = (phrase \land CAT : np \land DTR1 : word \land DTR1 : CAT : n$  $\wedge$  DTR1: PHON : Clinton  $\wedge$  DTR1: AGR : sg  $\wedge$  DTR2 : word  $\wedge$  DTR2: CAT :  $n \wedge$  DTR2: PHON : talks  $\wedge$  DTR2 :  $Z_2$ )  $\&\, \texttt{sign}(Z_2)$ r, c7,  $X = (phrase \land CAT : np \land DTR1 : word \land DTR1 : CAT : n$  $\wedge$  DTR1: PHON : Clinton  $\wedge$  DTR1: AGR : sg  $\wedge$  DTR2 : word  $\wedge$  DTR2: CAT :  $n \wedge$  DTR2: PHON : talks  $\wedge$  DTR2 :  $Z_2$ )

 $\& \operatorname{word}(Z_2)$ 

r, c

5,  $X = (phrase \land CAT : np \land DTR1 : word \land DTR1 : CAT : n \land DTR1 : PHON : Clinton \land DTR1 : AGR : sg \land DTR2 : word \land DTR2 : CAT : n \land DTR2 : PHON : talks \land DTR2 : AGR : pl)$ 

Figure 2.9: A derivation of  $[Clinton_N \ talks_N]_{NP}$ 

Furthermore, we reported the central formal details of feature-based CLGs and presented a simple linguistic grammar which will be used as a running example in the following chapters.

Proof trees or parses in constraint-based NLP can be quite complex even if simple grammars are used to analyze two-word phrases as in the example given above. Clearly, for complex grammars and phrases of reasonable length, structural ambiguity in constraint-based NLP is a severe problem. The task of the next two chapters is to provide a rigorous mathematical foundation of ambiguity resolution in constraint-based NLP.

# Chapter 3

# Quantitative CLP: Quantitative Inference with Subjective Weights and its Formal Semantics

In this chapter we present a novel framework for quantitative inference with subjective weights for CLP. We show soundness and completeness of the quantitative system with respect to a simple and intuitive formal semantics. We illustrate these concepts with a simple quantitative CLG and show how pruning techniques can be used to guide the search for the highest weighted analysis in such quantitative systems.

This chapter is based upon work previously published in Riezler (1996).

# 3.1 Introduction and Overview

Quantitative frameworks have been presented as extensions of both logic programming and constraint-based grammars. For the area of logic programming, a system of quantitative deduction which is sound and complete with respect to a related fixpoint semantics was introduced firstly by van Emden (1986). Like this seminal approach, most of the subsequent work on quantitative extensions of logic programming has concentrated on theoretical issues such as questions of the expressivity of systems for quantitative logic programming, or issues of the correctness of the connection of model-theory, fixpoint-theory, and proof-theory for such systems. However, none on these approaches seemed to have a specific application in mind.

On the contrary, quantitative extensions of constraint-based grammars have mainly been motivated by practical considerations. Most approaches in this area come as numerical extensions of the parsing strategy of existing constraint-based frameworks. However, even if for such systems the formal foundation of the underlying framework may be clear enough, none of these approaches comes with a well-defined semantics for its quantitative extension. That is, such quantitative extensions have to be seen as extralogical extensions of, e.g., the deduction scheme of the underlying CLP framework, and are not related to the model-theoretic counterpart of this operational semantics.

This is clearly an undesirable state of affairs. Rather, in the same way as CLGs provide a model-theoretic characterization of linguistic objects coupled with an operational parsing system, one would like to relate a quantitative deduction system to a quantitative modeltheory in a sound and complete way. The aim of this chapter is to present a sound and complete system of quantitative CLP which satisfies the following conditions. It should

- generally be applicable to CLP over arbitrary constraint languages,
- provide a precise, but yet simple formal semantics for quantitative CLP deduction,
- from the outset be designed with a specific application in mind, in our case, with respect to efficient ambiguity resolution in CLGs.

The first point means that in quantitative CLP one should not have to bother about the peculiarities of the constraint languages embedded into the CLP scheme. Rather, the quantitative extension should work in the same way for every constraint logic program irrespective of the embedded constraint language. For the NLP application, this means that for arbitrary constraint-based grammars a quantitative extension should be obtainable from the CLG resulting from an embedding of the grammar constraint language into a CLP scheme.

The second point addresses the tradeoff between the expressive power of the quantitative system and the intuitivity and simplicity of its semantics. That is, since the aim of a formal semantics is to provide a precise unambiguous way to specify the meaning of all aspects of an operational system at the design and implementation stage, it is justified only by its understandability and applicability. Our approach respects these ideas of simplicity and elegance by using the simple concepts of fuzzy set algebra as a basis for a formal semantics for quantitative CLP.

The third point, which refers to the intended application of ambiguity resolution and bestparse search in CLGs, is realized in quantitative CLP by stating the proof theory of quantitative CLP in terms of min/max trees, which in turn enables strategies such as alpha/betapruning to be used for efficient searching for best parses in CLGs.

Clearly, generalizations of this specific choice of design for quantitative CLP should be straightforward. However, they will not made explicit in the following chapters.

This chapter is organized as follows. Sect. 3.2 discusses previous work on quantitative logic programming and quantitative extensions of constraint-based grammars.

Sect. 3.3 introduces the concept of a quantitative definite clause specification, i.e., a quantitative constraint logic program.

Sect. 3.4 introduces the declarative semantics of quantitative definite clause specifications, i.e., a model-theoretic semantics based on concepts of fuzzy set algebra and a fixpoint semantics obtained by a minimal models in this model-theory.

Sect. 3.5 presents the operational semantics of quantitative CLP. That is, based on the concepts of quantitative derivation trees and quantitative proof trees, soundness and completeness of quantitative deduction in CLP is proven.

Sect. 3.6 exemplifies these concepts with a quantitative feature-based CLG, and shows how the search technique of alpha/beta-pruning can be applied to quantitative CLGs.

# 3.2 Previous Work

For the area of logic programming, van Emden (1986) presented in a seminal paper a quantitative deduction scheme and a fixpoint semantics for sets of numerically annotated Horn clauses. The aim of this paper was to enable the expression of a continuum of uncertainties between the usual two truth values in quantitative logic programs. The semantics of such quantitative logic programs is based upon concepts of fuzzy set algebra, and crucially deals with the truth-functional propagation of weights across conventional definite clauses. Van Emden's approach initialized research into a now extensively studied area of quantitative logic programming. For example, annotated logic programming (Subrahmanian (1987). Kifer and Subrahmanian (1992)) extends the expressive power of quantitative rule sets by allowing variables and evaluable function terms as annotations. Furthermore, in annotated logic programs, annotations can be attached to atoms and their conjunctions or disjunctions, and such programs are interpreted in powerful frameworks of lattice-theoretic semantics. Depending on different understandings of annotations, further extensions of van Emden (1986)'s and Subrahmanian (1987)'s approaches have been presented. Among those are approaches to possibilistic logic programming based on subjective necessity values (Dubois, Lang, and Prade 1991), probabilistic logic programming based on intervals of subjective probabilistic truth values (see, e.g., Ng and Subrahmanian (1992), Ng and Subrahmanian (1993)), or probabilistic deductive databases based on subjective confidence levels coming as intervals of belief and doubt (see, e.g., Lakshmanan and Sadri (1994), Lakshmanan and Sadri (1997)).

Quantitative extensions of constraint-based grammars have mainly been motivated by practical considerations. For example, Douglas and Dale (1992) presented an approach to robust parsing in PATR systems where according to a subjective value of necessity/optionality of constraints, constraint violations are allowed, and so robustness is introduced into the formalism. Kim (1994) presented an approach to best-first chart parsing with PATR grammars.

In this approach, atomic values of feature structures are annotated with subjective weights, and a weight combination scheme is defined for feature structure unification. The search space in best-first parsing then is restricted by a treshold below which completed and predicted feature structures are discarded. Erbach (1993a), Erbach (1993b), or Erbach (1998) introduced a model of preference for the CUF system, which is generalizable to the CLP scheme of Höhfeld and Smolka (1988), and which is used, among others, for tasks such as best-first parsing for ambiguity resolution and self-monitored generation. In Erbach's model, definite clauses as a whole are annotated with subjective preference values. Such preference values are combined in the resolution process by calculating the preference value of a clause consequent as the product of the preference value of the clause and the preference values of the antecedent predicates, which are additionally weighted to add up to 1.

The aim of our approach is to combine the mathematical exactness of the logicprogramming approaches with the practical applicability of the quantitative-grammar approaches. We will build our framework of quantitative CLP on ideas developed in the simple and elegant framework of van Emden (1986). This means that we restrict our attention to numerical weights attached to CLP clauses as a whole, and use the simple concepts of fuzzy set algebra to provide the basis for an intuitive formal semantics for quantitative CLP. Furthermore, we employ a min/max scheme for rule application which enables strategies such as alpha/beta pruning to be used for efficient searching. Clearly, our approach improves upon van Emden's approach by not being restricted to Horn clauses or to finite derivations. Moreover, it enables the application of quantitative search strategies to constraint-based grammars in a formally well-defined way.

# 3.3 Syntax of Quantitative CLP

Building upon the CLP scheme of Höhfeld and Smolka (1988) reported in Chap. 2, we can define the syntax of a quantitative definite clause specification  $\mathcal{P}_F$  very quickly. The following definitions are made with respect to implicit constraint languages  $\mathcal{L}$  and  $\mathcal{R}(\mathcal{L})$ . A definite clause specification  $\mathcal{P}$  in  $\mathcal{R}(\mathcal{L})$  then can be extended to a quantitative definite clause specification  $\mathcal{P}_F$  in  $\mathcal{R}(\mathcal{L})$  simply by adding numerical factors to program clauses.

**Definition 3.1** ( $\mathcal{P}_F$ ). A quantitative definite clause specification  $\mathcal{P}_F$  in  $\mathcal{R}(\mathcal{L})$  is a finite set of quantitative formulae, called quantitative definite clauses, of the form

$$\phi \& B_1 \& \dots \& B_n f \to A,$$

where  $A, B_1, \ldots, B_n$  are  $\mathcal{R}(\mathcal{L})$  -atoms,  $\phi$  is an  $\mathcal{L}$  -constraint,  $n \ge 0$ ,  $f \in (0, 1]$ . We may write a quantitative formula also as  $A \leftarrow_f \phi \& B_1 \& \ldots \& B_n$ .

These factors (the f in Definition 3.1) should be thought of as abstract weights which receive a concrete interpretation in specific instantiations of  $\mathcal{P}_F$ .

In the following the notation  $\mathcal{R}(\mathcal{L})$  will be used more generally to notate relationally extended constraint languages which possibly include quantitative formulae of the above form.

# 3.4 Declarative Semantics of Quantitative CLP

#### 3.4.1 Fuzzy Set Algebra and Model-Theoretic Semantics

To obtain a formal semantics for  $\mathcal{P}_F$ , first we have to introduce an appropriate quantitative measure into the set-theoretic specification of  $\mathcal{R}(\mathcal{L})$ -interpretations. One possibility to obtain quantitative  $\mathcal{R}(\mathcal{L})$ -interpretations is to base the set algebra of  $\mathcal{R}(\mathcal{L})$ -interpretations on the simple and well-defined concepts of fuzzy set algebra (see Zadeh (1965)).

Relying on Höhfeld and Smolka's specification of base equivalent  $\mathcal{R}(\mathcal{L})$ -interpretations, i.e.,  $\mathcal{R}(\mathcal{L})$ -interpretations extending the same  $\mathcal{L}$ -interpretation, in terms of the denotations of the relation symbols in these interpretations, we can "fuzzify" such interpretations by regarding the denotations of their relation symbols as fuzzy subsets of the set of tuples in the common domain.

Given constraint languages  $\mathcal{L}$  and  $\mathcal{R}(\mathcal{L})$ , we interpret each n-ary relation symbol  $r \in \mathcal{R}$ as a fuzzy subset of  $\mathcal{D}^n$ , for each  $\mathcal{R}(\mathcal{L})$ -interpretation  $\mathcal{A}$  with domain  $\mathcal{D}$ . That is, we identify the denotation of r under  $\mathcal{A}$  with a total function

$$\mu(\_; r^{\mathcal{A}}) : \mathcal{D}^n \to [0, 1],$$

which can be thought of as an abstract membership function. Such membership functions are generalized characteristic functions, and classical set membership is coded in this context by characteristic functions taking only 0 and 1 as values.

Next, we have to give a model-theoretic characterization of quantitative definite clauses. Clearly, any monotonous mapping could be used for the model-theoretic specification of the interaction of weights in quantitative definite clauses and accordingly for the calculation of weights in the proof-theory of quantitative CLP. For concreteness, we will instantiate such a mapping to the specific case of Definition 3.2 resembling van Emden (1986)'s mode of rule application. This will allow us to state the proof-theory of quantitative CLP in terms of min/max trees which in turn enables strategies such as alpha/beta pruning to be used for efficient searching. Such a quantitative CLP scheme improves upon several shortcomings of van Emden (1986)'s system, e.g. our quantitative CLP scheme clearly is not restricted to ground instances of Horn theories, and the soundness and completeness results we will present are not restricted to finite derivations. However, the choice of the mode of rule application made is not crucial for the substantial claims of this paper, and generalizations of this particular combination mode to specific applications should be straightforward, but are beyond the scope of this thesis.

The following definition of model corresponds to the definition of model in classical logic when considering only clauses with f = 1 and mappings  $\mathcal{D}^n \to \{0, 1\}$ .

**Definition 3.2 (Model).** An  $\mathcal{R}(\mathcal{L})$ -interpretation  $\mathcal{A}$  extending some  $\mathcal{L}$ -interpretation  $\mathcal{I}$  is a model of a quantitative definite clause specification  $\mathcal{P}_F$  iff for each  $\alpha \in \mathsf{ASS}$ , for each quantitative formula  $r(\vec{x}) \leftarrow_f \phi \& q_1(\vec{x}_1) \& \ldots \& q_k(\vec{x}_k)$  in  $\mathcal{P}_F$  holds:

If 
$$\alpha \in \llbracket \phi \rrbracket^{\mathcal{I}}$$
, then  $\mu(\alpha(\vec{x}); r^{\mathcal{A}}) \ge f \times \min\{\mu(\alpha(\vec{x}_i); q_i^{\mathcal{A}}) | 1 \le j \le k\}$ .

In terms of membership degrees, this definition of model can be paraphrased as follows: If the antecedent constraint is satisfiable, then the membership degrees of the denotations of the consequent atom must not be less than f times the membership degrees of the denotations of the antecedent atom. A truth-functional view could be obtained by considering membership degrees as truth degrees of atoms under variable assignments. From the viewpoint of such a truth-functional propagation of weights across definite clauses, a clause contributes to the consequent a truth value which is f times the truth value of the antecedent.

Note that the notation of an  $\mathcal{R}(\mathcal{L})$ -interpretation  $\mathcal{A}$  will be used more generally to include interpretations of quantitative formulae.  $\mathcal{R}(\mathcal{L})$ -solutions of a quantitative formula are defined as  $[\![r(\vec{x}) \leftarrow_f \phi \& q_1(\vec{x}_1) \& \dots \& q_k(\vec{x}_k)]\!]^{\mathcal{A}} = \{ \alpha \in \mathsf{ASS} \mid \text{ If } \alpha \in [\![\phi]\!]^{\mathcal{I}}, \text{ then } \mu(\alpha(\vec{x}); r^{\mathcal{A}}) \geq f \times \min\{\mu(\alpha(\vec{x}_j); q_j^{\mathcal{A}}) \mid 1 \leq j \leq k\} \}.$ 

Based on the above definition of model, the concept of logical consequence can be defined as usual.

**Definition 3.3 (Logical consequence).** A quantitative formula  $r(\vec{x}) \leftarrow_f \phi$  is a logical consequence of a quantitative definite clause specification  $\mathcal{P}_F$  iff for each  $\mathcal{R}(\mathcal{L})$ -interpretation  $\mathcal{A}$ ,  $\mathcal{A}$  is a model of  $\mathcal{P}_F$  implies that  $\mathcal{A}$  is a model of  $\{r(\vec{x}) \leftarrow_f \phi\}$ .

Furthermore, we have that the fact that  $r(\vec{x}) \leftarrow_f \phi$  is a logical consequence of  $\mathcal{P}_F$  implies that  $r(\vec{x}) \leftarrow_{f'} \phi$  is a logical consequence of  $\mathcal{P}_F$  for every  $f' \leq f$ .

A goal G is defined similar to the non-quantitative case as a (possibly empty) conjunction of  $\mathcal{R}(\mathcal{L})$ -atoms and  $\mathcal{L}$ -constraints. We can, without loss of generality, restrict goals to be of the form  $r(\vec{x}) \& \phi$ , i.e., a (possibly empty) conjunction of a single relational atom  $r(\vec{x})$ and an  $\mathcal{L}$ -constraint  $\phi$ . This can be done since for each goal  $G = r_1(\vec{x}_1) \& \ldots \& r_k(\vec{x}_k) \& \phi$ which contains more than one relational atom, we can complete the program with a new clause  $C = r(\vec{x}_1, \ldots, \vec{x}_k) \leftarrow_1 r_1(\vec{x}_1) \& \ldots \& r_k(\vec{x}_k) \& \phi$ , with G as antecedent and a new predicate, which takes all variables in G as arguments, as consequent. Submitting the new predicate  $r(\vec{x}_1, \ldots, \vec{x}_k)$  as query yields the same results as would be obtained when querying with the compound goal G.

Given some program  $\mathcal{P}_F$  and some goal G, a quantitative  $\mathcal{P}_F$ -answer  $\varphi$  of G is defined as a satisfiable  $\mathcal{L}$ -constraint  $\varphi$  s.t.  $\varphi_f \to G$  is a logical consequence of  $\mathcal{P}_F$ . A quantitative formula  $\varphi_f \to r(\vec{x}) \& \phi$  is defined to be a logical consequence of  $\mathcal{P}_F$  iff every model of  $\mathcal{P}_F$  is a model of  $\{\varphi_f \to r(\vec{x}) \& \phi\}$ . An  $\mathcal{R}(\mathcal{L})$ -interpretation  $\mathcal{A}$  is a model of  $\{\varphi_f \to r(\vec{x}) \& \phi\}$  iff  $[\![\varphi]\!]^{\mathcal{A}} \subseteq [\![\phi]\!]^{\mathcal{A}}$  and  $\mathcal{A}$  is a model of  $\{r(\vec{x}) \leftarrow_f \varphi\}$ .

Next we have to associate a complete lattice of interpretations with quantitative definite clause specifications.

Adopting Zadeh's definitions for set operations, we can define a partial ordering on the set of base equivalent  $\mathcal{R}(\mathcal{L})$ -interpretations. This is done by defining set operations on these interpretations with reference to set operations on the denotations of relation symbols in these interpretations. We get for all base equivalent  $\mathcal{R}(\mathcal{L})$ -interpretations  $\mathcal{A}$ ,  $\mathcal{A}'$ :

- $\mathcal{A} \subseteq \mathcal{A}'$  iff for each n-ary relation symbol  $r \in \mathcal{R}$ , for each  $\alpha \in ASS$ , for each  $\vec{x} \in VAR^n$ :  $\mu(\alpha(\vec{x}); r^{\mathcal{A}}) \leq \mu(\alpha(\vec{x}); r^{\mathcal{A}'}),$
- $\mathcal{A} = \bigcup X$  iff for each n-ary relation symbol  $r \in \mathcal{R}$ , for each  $\alpha \in \mathsf{ASS}$ , for each  $\vec{x} \in \mathsf{VAR}^n$ :  $\mu(\alpha(\vec{x}); r^{\mathcal{A}}) = \sup\{\mu(\alpha(\vec{x}); r^{\mathcal{A}'}) | \mathcal{A}' \in X\},\$
- $\mathcal{A} = \bigcap X$  iff for each n-ary relation symbol  $r \in \mathcal{R}$ , for each  $\alpha \in \mathsf{ASS}$ , for each  $\vec{x} \in \mathsf{VAR}^n$ :  $\mu(\alpha(\vec{x}); r^{\mathcal{A}}) = \inf\{\mu(\alpha(\vec{x}); r^{\mathcal{A}'}) | \mathcal{A}' \in X\}.$

Note that we define furthermore  $\sup \emptyset = 0$ ,  $\inf \emptyset = 1$ . Clearly, the set of all base equivalent  $\mathcal{R}(\mathcal{L})$ -interpretations is a complete lattice under the partial ordering of set inclusion. The supremum is given by the union, and the infimum by the intersection, for any set of base-equivalent  $\mathcal{R}(\mathcal{L})$ -interpretations. The top element is the  $\mathcal{R}(\mathcal{L})$ -interpretation  $\mathcal{A}^{\top}$  such that for each  $r \in \mathcal{R}$ , for each  $\vec{u} \in \mathcal{D}^{\mathsf{Ar}(r)}$ :  $\mu(\vec{u}; r^{\mathcal{A}^{\top}}) = 1$ , and the bottom element is the  $\mathcal{R}(\mathcal{L})$ -interpretation  $\mathcal{A}^{\perp}$  such that for each  $r \in \mathcal{R}$ , for each  $\vec{u} \in \mathcal{D}^{\mathsf{Ar}(r)}$ :  $\mu(\vec{u}; r^{\mathcal{A}^{\perp}}) = 0$ .

#### 3.4.2 Minimal Model Semantics

Based upon the definition of a complete lattice of  $\mathcal{R}(\mathcal{L})$ -interpretations of a quantitative definite clause specification  $\mathcal{P}_F$ , we can state the following equations, which link the declarative and operational semantics of  $\mathcal{P}_F$ . These equations define the notion of a  $\mathcal{P}_F$ -chain, which will be crucial for the construction of minimal models for  $\mathcal{P}_F$ . Similar to the non-quantitative case, these equations are based on the respective definition of model, and take for the quantitative case the following form.

**Definition 3.4.** Let  $\mathcal{P}_F$  be a quantitative definite clause specification in  $\mathcal{R}(\mathcal{L})$ ,  $\mathcal{I}$  be an  $\mathcal{L}$ interpretation. Then the countably infinite sequence  $\langle \mathcal{A}_0, \mathcal{A}_1, \mathcal{A}_2, \ldots \rangle$  of  $\mathcal{R}(\mathcal{L})$ -interpretations
extending  $\mathcal{I}$  is a  $\mathcal{P}_F$ -chain iff for each n-ary relation symbol  $r \in \mathcal{R}$ , for each  $\alpha \in ASS$ , for
each  $\vec{x} \in VAR^n$ :

 $\mu(\alpha(\vec{x}); r^{\mathcal{A}_0}) := 0,$ 

 $\mu(\alpha(\vec{x}); r^{\mathcal{A}_{i+1}}) := \max\{f \times \min\{\mu(\alpha(\vec{x}_j); q_j^{\mathcal{A}_i}) | 1 \le j \le n\} \mid \text{there is a variant } r(\vec{x}) \leftarrow_f \phi \& q_1(\vec{x}_1) \& \dots \& q_n(\vec{x}_n) \text{ of a clause in } \mathcal{P}_F \text{ and } \alpha \in [\![\phi]\!]^{\mathcal{A}_i} \}.$ 

Before turning to the construction of minimal models, we have to prove the following useful lemma (see van Emden (1986), Lemmata 2.10', 2.11'). Lemma 3.1 assures that for each tuple of objects in the denotation of a relation symbol under a minimal model, there is a corresponding finite step in the  $\mathcal{P}_F$ -chain which introduces these objects into the minimal model denotation.

**Lemma 3.1.** For each  $\mathcal{P}_F$ , for each  $\mathcal{P}_F$ -chain  $\langle \mathcal{A}_0, \mathcal{A}_1, \mathcal{A}_2, \ldots \rangle$ , for each k-ary relation symbol  $r \in \mathcal{R}$ , for each  $\alpha \in ASS$ , for each  $\vec{x} \in VAR^k$ , there exists some  $n \in \mathbb{N}$  s.t.  $\mu(\alpha(\vec{x}); r^{\bigcup_{i \geq 0} \mathcal{A}_i}) = \mu(\alpha(\vec{x}); r^{\mathcal{A}_n}).$ 

*Proof.* We have to show that the supremum  $v = \sup\{\mu(\alpha(\vec{x}); r^{\mathcal{A}_i}) | i \ge 0\}$  can be attained for some  $n \in \mathbb{N}$ .

- v = 0: For v = 0, we have n = 0.
- v > 0: For v > 0, we have to show that for any real  $\epsilon$ ,  $0 < \epsilon < v$ , the set  $\{\mu(\alpha(\vec{x}); r^{\mathcal{A}_i}) | i \geq 0 \text{ and } \mu(\alpha(\vec{x}); r^{\mathcal{A}_i}) \geq \epsilon\}$  is finite.

Let F be the finite set of real numbers of factors of clauses in  $\mathcal{P}_F$ , m be the greatest element in F s.t. m < 1 and let q be the smallest integer s.t.  $m^q < \epsilon$ .

Then, since each real number  $\mu(\alpha(\vec{x}); r^{\mathcal{A}_i})$  is a product of a sequence of elements of F, the number of different products  $\geq \epsilon$  is not greater than  $|F|^q$ , the permutation of |F| different things taken q at a time with repetitions, and thus finite.

Hence, the supremum is the maximum attained for some  $n \in \mathbb{N}$ .

Now we can obtain minimal model properties for quantitative definite clause specifications similar to those for the non-quantitative programs of Höhfeld and Smolka (1988). Based on the constructive definition of a  $\mathcal{P}_F$ -chain of  $\mathcal{R}(\mathcal{L})$ -interpretations extending an  $\mathcal{L}$ -interpretation  $\mathcal{I}$ , an  $\mathcal{R}(\mathcal{L})$ -interpretation  $\mathcal{A}$  is obtainable as the  $\mathcal{R}(\mathcal{L})$ -interpretation which is both a model of  $\mathcal{P}_F$  and minimal with respect to the lattice of base equivalent  $\mathcal{R}(\mathcal{L})$ -interpretations extending  $\mathcal{I}$ . Theorem 3.2 states that we can construct a minimal model  $\mathcal{A}$  of  $\mathcal{P}_F$  for each quantitative definite clause specification  $\mathcal{P}_F$  in the extension of an arbitrary constraint language  $\mathcal{L}$  and for each  $\mathcal{L}$ -interpretation. This means that—due to the definiteness of  $\mathcal{P}_F$ —we can restrict our attention to a minimal model semantics of  $\mathcal{P}_F$ .

**Theorem 3.2 (Definiteness).** For each  $\mathcal{L}$ -interpretation  $\mathcal{I}$ , for each quantitative definite clause specification  $\mathcal{P}_F$  in  $\mathcal{R}(\mathcal{L})$ , for each  $\mathcal{P}_F$ -chain  $\langle \mathcal{A}_0, \mathcal{A}_1, \mathcal{A}_2, \ldots \rangle$  of  $\mathcal{R}(\mathcal{L})$ -interpretations extending some  $\mathcal{L}$ -interpretation  $\mathcal{I}$ :

- (i)  $\mathcal{A}_0 \subseteq \mathcal{A}_1 \subseteq \ldots$ ,
- (ii) the union  $\mathcal{A} := \bigcup_{i>0} \mathcal{A}_i$  is a model of  $\mathcal{P}_F$  extending  $\mathcal{I}$ ,
- (iii)  $\mathcal{A}$  is the minimal model of  $\mathcal{P}_F$  extending  $\mathcal{I}$ .

*Proof.* (i) We have to show that  $\mathcal{A}_i \subseteq \mathcal{A}_{i+1}$ . We prove by induction on i showing for each constraint language  $\mathcal{L}$ , for each quantitative definite clause specification  $\mathcal{P}_F$  in  $\mathcal{R}(\mathcal{L})$ , for each  $\mathcal{L}$ -interpretation  $\mathcal{I}$ , for each  $\mathcal{P}_F$ -chain  $\langle \mathcal{A}_0, \mathcal{A}_1, \mathcal{A}_2, \ldots \rangle$  of  $\mathcal{R}(\mathcal{L})$ -interpretations extending some  $\mathcal{L}$ -interpretation  $\mathcal{I}$ , for each n-ary relation symbol  $r \in \mathcal{R}$ , for each  $\alpha \in ASS$ , for each  $\vec{x} \in VAR^n$ , for each  $i \in \mathbb{N}$ :  $\mu(\alpha(\vec{x}); r^{\mathcal{A}_i}) \leq \mu(\alpha(\vec{x}); r^{\mathcal{A}_{i+1}})$ .

Base:  $\mu(\alpha(\vec{x}); r^{\mathcal{A}_0}) = 0 \le \mu(\alpha(\vec{x}); r^{\mathcal{A}_1}).$ 

Hypothesis: Suppose  $\mu(\alpha(\vec{x}); r^{\mathcal{A}_{n\perp 1}}) \leq \mu(\alpha(\vec{x}); r^{\mathcal{A}_n}).$ 

Step:  $\mu(\alpha(\vec{x}); r^{\mathcal{A}_n}) = v > 0$ 

- $\implies \text{ there exists a variant } r(\vec{x}) \leftarrow_f \phi \& q_1(\vec{x}_1) \& \dots \& q_k(\vec{x}_k) \text{ of a clause in } \mathcal{P}_F \text{ s.t. } v = f \times \min\{\mu(\alpha(\vec{x}_1); q_1^{\mathcal{A}_{n\perp 1}}), \dots, \mu(\alpha(\vec{x}_k); q_k^{\mathcal{A}_{n\perp 1}})\} \text{ and } \alpha \in \llbracket\phi\rrbracket^{\mathcal{A}_{n\perp 1}}, \text{ by Definition} 3.4$
- $\implies \mu(\alpha(\vec{x}_1); q_1^{\mathcal{A}_n}) \ge \mu(\alpha(\vec{x}_1); q_1^{\mathcal{A}_{n\perp 1}}), \dots, \mu(\alpha(\vec{x}_k); q_k^{\mathcal{A}_n}) \ge \mu(\alpha(\vec{x}_k); q_k^{\mathcal{A}_{n\perp 1}}) \text{ and } \alpha \in \llbracket \phi \rrbracket^{\mathcal{A}_n}, \text{ by the hypothesis}$
- $\implies \mu(\alpha(\vec{x}); r^{\mathcal{A}_{n+1}}) \ge v$ , by definition of  $\mu(\alpha(\vec{x}); r^{\mathcal{A}_{i+1}})$

$$\implies \mu(\alpha(\vec{x}); r^{\mathcal{A}_n}) \le \mu(\alpha(\vec{x}); r^{\mathcal{A}_{n+1}})$$

For v = 0 it follows immediately that  $\mu(\alpha(\vec{x}); r^{\mathcal{A}_n}) \leq \mu(\alpha(\vec{x}); r^{\mathcal{A}_{n+1}}).$ 

Claim (i) follows by arithmetic induction.

(ii) We have to show that  $\mathcal{A} := \bigcup_{i\geq 0} \mathcal{A}_i$  is a model of  $\mathcal{P}_F$  extending  $\mathcal{I}$ . We prove that for each clause  $r(\vec{x}) \leftarrow_f \phi \& q_1(\vec{x}_1) \& \dots \& q_k(\vec{x}_k)$  in  $\mathcal{P}_F$ , for each  $\alpha \in \mathsf{ASS}$ : If  $\alpha \in \llbracket \phi \rrbracket^{\mathcal{A}}$ , then  $\mu(\alpha(\vec{x}); r^{\mathcal{A}}) \geq f \times \min\{\mu(\alpha(\vec{x}_j); q_j^{\mathcal{A}}) \mid 1 \leq j \leq k\}.$ 

- Note that since every  $\mathcal{A}_i$  is an  $\mathcal{R}(\mathcal{L})$ -interpretation extending  $\mathcal{I}$ ,  $\mathcal{A}$  is an  $\mathcal{R}(\mathcal{L})$ -interpretation extending  $\mathcal{I}$ .
- Now let  $r(\vec{x}) \leftarrow_f \phi \& q_1(\vec{x}_1) \& \dots \& q_k(\vec{x}_k)$  be a clause in  $\mathcal{P}_F$  s.t. for some  $\alpha \in \mathsf{ASS}$ :  $\alpha \in \llbracket \phi \rrbracket^{\mathcal{A}}$  and  $\mu(\alpha(\vec{x}_i); q_i^{\mathcal{A}}) = \min\{\mu(\alpha(\vec{x}_i); q_i^{\mathcal{A}}) | 1 \le j \le k\} = v.$

Then there exists some 
$$n \in \mathbb{N}$$
 s.t.  $v = \mu(\alpha(\vec{x}_i); q_i^{\mathcal{A}_n}) = \min\{\mu(\alpha(\vec{x}_j); q_j^{\mathcal{A}_n}) | 1 \le j \le k\}$ , by  
Lemma 3.1 and since for all  $j$  s.t.  $1 \le j \le k : \mu(\alpha(\vec{x}_j); q_j^{\mathcal{A}}) = \sup\{\mu(\alpha(\vec{x}_j); q_j^{\mathcal{A}_i}) | i \ge 0\}$ 

$$\implies \mu(\alpha(\vec{x}); r^{\mathcal{A}_{n+1}}) \ge f \times v, \text{ by Definition 3.4}$$
$$\implies \mu(\alpha(\vec{x}); r^{\mathcal{A}}) \ge \mu(\alpha(\vec{x}); r^{\mathcal{A}_{n+1}}), \text{ since } \mu(\alpha(\vec{x}); r^{\mathcal{A}}) = \sup\{\mu(\alpha(\vec{x}); r^{\mathcal{A}_i}) | i \ge 0\}$$
$$\implies \mu(\alpha(\vec{x}); r^{\mathcal{A}}) \ge f \times \min\{\mu(\alpha(\vec{x}_j); q_j^{\mathcal{A}}) | 1 \le j \le k\}.$$

This completes the proof for claim (ii).

(iii) We have to show that  $\mathcal{A}$  is the minimal model of  $\mathcal{P}_F$  extending  $\mathcal{I}$ . We prove for every base equivalent model  $\mathcal{B}$  of  $\mathcal{P}_F : \mathcal{A}_i \subseteq \mathcal{B}$ , which gives  $\mathcal{A} \subseteq \mathcal{B}$ , by induction on *i* showing for each constraint language  $\mathcal{L}$ , for each quantitative definite clause specification  $\mathcal{P}_F$  in  $\mathcal{R}(\mathcal{L})$ , for each  $\mathcal{L}$ -interpretation  $\mathcal{I}$ , for each  $\mathcal{P}_F$ -chain  $\langle \mathcal{A}_0, \mathcal{A}_1, \mathcal{A}_2, \ldots \rangle$  of  $\mathcal{R}(\mathcal{L})$ -interpretations extending some  $\mathcal{L}$ -interpretation  $\mathcal{I}$ , for each n-ary relation symbol  $r \in \mathcal{R}$ , for each  $\alpha \in ASS$ , for each  $\vec{x} \in VAR^n$ , for each  $i \in \mathbb{N}$ :  $\mu(\alpha(\vec{x}); r^{\mathcal{A}_i}) \leq \mu(\alpha(\vec{x}); r^{\mathcal{B}})$ .

Base:  $\mu(\alpha(\vec{x}); r^{\mathcal{A}_0}) = 0 \le \mu(\alpha(\vec{x}); r^{\mathcal{B}}).$ 

Hypothesis: Suppose  $\mu(\alpha(\vec{x}); r^{\mathcal{A}_{n\perp 1}}) \leq \mu(\alpha(\vec{x}); r^{\mathcal{B}}).$ 

Step:  $\mu(\alpha(\vec{x}); r^{\mathcal{A}_n}) = v > 0$ 

- $\implies \text{there exists a variant } r(\vec{x}) \leftarrow_f \phi \& q_1(\vec{x}_1) \& \dots \& q_k(\vec{x}_k) \text{ of a clause in } \mathcal{P}_F \text{ s.t. } v = f \times \min\{\mu(\alpha(\vec{x}_1); q_1^{\mathcal{A}_{n\perp 1}}), \dots, \mu(\alpha(\vec{x}_k); q_k^{\mathcal{A}_{n\perp 1}})\} \text{ and } \alpha \in \llbracket\phi\rrbracket^{\mathcal{A}_{n\perp 1}}, \text{ by Definition} 3.4$
- $\implies \mu(\alpha(\vec{x}_1); q_1^{\mathcal{B}}) \geq \mu(\alpha(\vec{x}_1); q_1^{\mathcal{A}_{n\perp 1}}), \dots, \mu(\alpha(\vec{x}_k); q_k^{\mathcal{B}}) \geq \mu(\alpha(\vec{x}_k); q_k^{\mathcal{A}_{n\perp 1}}) \text{ and } \alpha \in \llbracket \phi \rrbracket^{\mathcal{B}}, \text{ by the hypothesis}$

$$\implies \mu(\alpha(\vec{x}); r^{\mathcal{B}}) \geq v$$
, since  $\mathcal{B}$  is a model of  $\mathcal{P}_F$ 

 $\implies \mu(\alpha(\vec{x}); r^{\mathcal{A}_n}) \le \mu(\alpha(\vec{x}); r^{\mathcal{B}}).$ 

For v = 0 it follows immediately that  $\mu(\alpha(\vec{x}); r^{\mathcal{A}_n}) \leq \mu(\alpha(\vec{x}); r^{\mathcal{B}}).$ 

Claim (iii) follows by arithmetic induction.

The following proposition allows us to link the declarative description of the desired output from  $\mathcal{P}_F$  and a goal, i.e., a quantitative  $\mathcal{P}_F$ -answer, to the minimal model semantics of  $\mathcal{P}_F$ . That is, Proposition 3.3 shows that quantitative  $\mathcal{P}_F$ -answers are completely characterized by minimal models of  $\mathcal{P}_F$ . Similar to the non-quantitative case, this is done for the quantitative case by connecting the concept of logical consequence with the concept of minimal model.

**Proposition 3.3.** Let  $\mathcal{P}_F$  be a quantitative definite clause specification in  $\mathcal{R}(\mathcal{L})$ ,  $\varphi$  be an  $\mathcal{L}$ constraint and G be a goal. Then  $\varphi_v \to G$  is a logical consequence of  $\mathcal{P}_F$  iff every minimal
model  $\mathcal{A}$  of  $\mathcal{P}_F$  is a model of  $\{\varphi_v \to G\}$ .

*Proof.* If: For each minimal model  $\mathcal{A}$  of  $\mathcal{P}_F : \mathcal{A}$  is a model of  $\{\varphi_v \to G\}$ 

 $\implies \text{for every model } \mathcal{B} \text{ of } \mathcal{P}_F \text{ base equivalent to some minimal model } \mathcal{A} \text{ of } \mathcal{P}_F : \mathcal{B} \text{ is a} \\ \text{model of } \{\varphi_v \to G\}, \text{ since } \mathcal{A} \subseteq \mathcal{B} \text{ by Theorem 3.2, (iii)} \\ \implies \varphi_v \to G \text{ is a logical consequence of } \mathcal{P}_F.$ 

Only if:  $\varphi_v \to G$  is a logical consequence of  $\mathcal{P}_F$ 

$$\implies$$
 every model of  $\mathcal{P}_F$  is a model of  $\{\varphi_v \to G\}$ , by Definition 3.3

$$\implies \mathcal{A} \text{ is a model of } \{\varphi_v \to G\}.$$

The following example illustrates the basic concepts of the declarative semantics of quantitative definite clause specifications. The program of Fig. 3.1 is a quantitative version of the program of Fig. 2.1. The factors attached to clauses 2 and 3 express a preference of the  $\mathcal{L}$ constraint X = a over the  $\mathcal{L}$ -constraint X = b in the definition of the predicate p. Predicate q is defined uniquely in clause 1 and gets assigned the factor 1.

1 
$$q(X) \leftarrow_1 p(X)$$
.  
2  $p(X) \leftarrow_{.7} X = a$   
3  $p(X) \leftarrow_{.5} X = b$ .

Figure 3.1: Quantitative constraint logic program

The construction of a minimal model for the program of Fig. 3.1 is shown in Fig. 3.2. For a variable assignment  $\alpha \in [X = a]^{\mathcal{I}}$ , the membership value of .7 of the object  $\langle \alpha(X) \rangle$  in the denotation of the predicate **p** (resp. **q**) under the minimal model  $\mathcal{A}$  is obtained in step 1 (resp. step 2) of the  $\mathcal{P}_F$ -chain construction. For a variable assignment  $\alpha \in [X = b]^{\mathcal{I}}$ , a membership degree of .5 is obtained in similar manner.

Clearly,  $\mathcal{A} = \bigcup_{i>0} \mathcal{A}_i$  is a minimal model of the quantitative program of Fig. 3.1.

# 3.5 Operational Semantics for Quantitative CLP

#### 3.5.1 Min/Max Trees and Quantitative Proof Trees

The proof procedure for quantitative CLP can be stated conveniently as a search of a tree, corresponding to the search of an SLD-and/or tree in conventional logic programming or to the search of a derivation tree as defined in Chap. 2 for CLP. The structure of such a tree exactly mirrors the construction of a minimal model and thus may be defined as a min/max tree. That is, according to the minimal model construction, which is based on the operations min and max, a min/max tree combines the standard left-right selection and depth-first search

$$\begin{split} \alpha \in \llbracket X &= a \rrbracket^{\mathcal{I}}: \\ \mu(\langle \alpha(X) \rangle \, ; \, \mathfrak{p}^{\mathcal{A}_{0}}) &= 0, \\ \mu(\langle \alpha(X) \rangle \, ; \, \mathfrak{p}^{\mathcal{A}_{1}}) &= \max\{.7 \times \min \emptyset\} = .7, \\ \mu(\langle \alpha(X) \rangle \, ; \, \mathfrak{p}^{\mathcal{A}_{2}}) &= \max\{.7 \times \min \emptyset\} = .7, \\ \vdots \\ \mu(\langle \alpha(X) \rangle \, ; \, \mathfrak{p}^{\bigcup_{i \geq 0} \mathcal{A}_{i}}) &= \sup\{0, .7, .7, \dots\} = .7, \end{split}$$

$$\mu(\langle \alpha(X) \rangle; \mathbf{q}^{\mathcal{A}_0}) = 0,$$
  

$$\mu(\langle \alpha(X) \rangle; \mathbf{q}^{\mathcal{A}_1}) = 0,$$
  

$$\mu(\langle \alpha(X) \rangle; \mathbf{q}^{\mathcal{A}_2}) = \max\{1 \times \min\{.7\}\} = .7,$$
  

$$\vdots$$
  

$$\mu(\langle \alpha(X) \rangle; \mathbf{q}^{\bigcup_{i \ge 0} \mathcal{A}_i}) = \sup\{0, 0, .7, \dots\} = .7.$$

$$\begin{split} &\alpha \in [\![X = b]\!]^{\mathcal{I}}: \\ &\mu(\langle \alpha(X) \rangle \, ; \, \mathbf{p}^{\mathcal{A}_0}) = 0, \\ &\mu(\langle \alpha(X) \rangle \, ; \, \mathbf{p}^{\mathcal{A}_1}) = \max\{.5 \times \min \emptyset\} = .5, \\ &\mu(\langle \alpha(X) \rangle \, ; \, \mathbf{p}^{\mathcal{A}_2}) = \max\{.5 \times \min \emptyset\} = .5, \\ &\vdots \\ &\mu(\langle \alpha(X) \rangle \, ; \, \mathbf{p}^{\bigcup_{i \ge 0} \mathcal{A}_i}) = \sup\{0, .5, .5, \dots\} = .5, \end{split}$$

$$\mu(\langle \alpha(X) \rangle; \mathbf{q}^{\mathcal{A}_0}) = 0,$$
  

$$\mu(\langle \alpha(X) \rangle; \mathbf{q}^{\mathcal{A}_1}) = 0,$$
  

$$\mu(\langle \alpha(X) \rangle; \mathbf{q}^{\mathcal{A}_2}) = \max\{1 \times \min\{.5\}\} = .5,$$
  

$$\vdots$$
  

$$\mu(\langle \alpha(X) \rangle; \mathbf{q}^{\bigcup_{i \ge 0} \mathcal{A}_i}) = \sup\{0, 0, .5, \dots\} = .5,$$

Figure 3.2:  $\mathcal{P}_F$  -chain for quantitative constraint logic program

with a min/max calculation of node-values. A relation node of a derivation tree corresponds in the quantitative case to a max-node, and a constraint node to a min-node. In contrast to derivation trees, in min/max trees the unique successor of a constraint node is split up into several successor nodes, one for each relational atom in the goal. This is necessary to calculate a minimum of node values at a min-node.

In the following we will assume implicit constraint languages  $\mathcal{L}$  and  $\mathcal{R}(\mathcal{L})$  and a given quantitative definite clause specification  $\mathcal{P}_F$  in  $\mathcal{R}(\mathcal{L})$ . Furthermore, V will denote the finite set of variables in the query and the V-solutions of a constraint  $\phi$  in an interpretation  $\mathcal{I}$  are defined as  $\llbracket \phi \rrbracket_V^{\mathcal{I}} := \{ \alpha |_V | \ \alpha \in \llbracket \phi \rrbracket^{\mathcal{I}} \}$  and  $\alpha |_V$  is the restriction of  $\alpha$  to V.

**Definition 3.5 (Min/max tree).** A min/max tree determined by a query  $G_1$  and a quantitative definite clause specification  $\mathcal{P}_F$  has to satisfy the following conditions:

- Each max-node is labeled by a goal. The value of each nonterminal max-node is the maximum of the values of its successors.
- Each min-node is labeled by a clause from  $\mathcal{P}_F$  and a goal. The value of each nonterminal min-node is  $f \times m$ , where f is the factor of the clause and m is the minimum of the values of its successors.
- The successors of every max-node are all min-nodes s.t. for every clause C with <sup>r</sup>→ resolvent G' obtained by C from goal G in a max-node, there is a min-node successor
  labeled by C and G'.
- The successors of every min-node are all max-nodes s.t. for every  $\mathcal{R}(\mathcal{L})$ -atom  $r(\vec{x})$  in goal  $G \& \phi \& \phi'$  in a min-node with  $\stackrel{c}{\longrightarrow}$ -resolvent  $G \& \phi''$ , there is a max-node successor labeled by  $r(\vec{x}) \& \phi''$ .
- The root node is a max-node labeled by  $G_1$ .
- A success node is a terminal max-node labeled by a satisfiable  $\mathcal{L}$  -constraint. The value of a success node is 1.
- A failure node is a terminal max-node which is not a success node. The value of a failure node is 0.

Similar to the non-quantitative case, a proof tree in the quantitative case is a subtree of a derivation tree. However, in a quantitative proof tree, each min-node takes all of the successors of the min-node of the min/max tree as its successors. Furthermore, to check the consistency of the constraint solving results in the min-node successors, an additional  $\stackrel{c}{\longrightarrow}$ -step has to be applied to the conjunction of all success nodes of a quantitative proof tree. This step yields a satisfiable  $\mathcal{L}$ -constraint, called answer constraint, if the conjunction of the  $\mathcal{L}$ -constraints in the success nodes is satisfiable.

**Definition 3.6 (Quantitative proof tree).** A quantitative proof tree for a goal  $G_1$  from quantitative definite clause specification  $\mathcal{P}_F$  is a subtree of a min/max supertree determined by  $G_1$  and  $\mathcal{P}_F$  and defined as follows:

- The root node of the proof tree is the root node of the supertree.
- A max-node of the proof tree is a max-node of the supertree and takes one of the successors of the supertree max-node as its successor.
- A min-node of the proof tree is a min-node of the supertree and takes all of the successors of the supertree max-node as its successors.
- All terminal nodes in the proof tree are success nodes  $\phi, \phi', \dots s.t. \phi \& \phi' \& \dots \xrightarrow{c} \varphi$ and  $\varphi$  is a satisfiable  $\mathcal{L}$ -constraint, called answer constraint.
- Values are assigned to proof tree nodes in the same way as to min/max tree nodes.

#### 3.5.2 Soundness and Completeness

To prove soundness and completeness of the generalized SLD-resolution proof procedure defined via min/max trees and quantitative proof trees, some further concepts have to be introduced.

Note that the definitions of renaming,  $\rho$ -variant, and variant carry over to the quantitative case without changes. Clearly, we have the property that a constraint language  $\mathcal{R}(\mathcal{L})$  containing quantitative definite clauses is closed under renaming if the underlying constraint language  $\mathcal{L}$  is closed under renaming. Furthermore, for each such generalized constraint language  $\mathcal{R}(\mathcal{L})$  which is closed under renaming, and for each  $\mathcal{R}(\mathcal{L})$ -interpretation  $\mathcal{A}$ , we have that  $\mathcal{A}$  is a model of an  $\mathcal{R}(\mathcal{L})$ -constraint iff  $\mathcal{A}$  is a model of each of its variants.

Next, we have to redefine a complexity measure for goal reduction for the quantitative case. This measure is crucial in proving termination of goal reduction and works by keying steps of the minimal model construction to steps of the goal reduction process.

• The complexity of a variable assignment  $\alpha$  for an atom  $r(\vec{x})$  in the minimal model  $\mathcal{A}$  s.t.  $\mu(\alpha(\vec{x}); r^{\mathcal{A}}) > 0$  is defined as

$$comp(\alpha, r(\vec{x}), \mathcal{A}) := \min\{i \mid \mu(\alpha(\vec{x}); r^{\mathcal{A}}) = \mu(\alpha(\vec{x}); r^{\mathcal{A}_i})\};$$

• The complexity of  $\alpha$  for goal  $G = r_1(\vec{x}_1)$  & ... &  $r_k(\vec{x}_k)$  &  $\phi$  in  $\mathcal{A}$  s.t.  $\alpha \in \llbracket \phi \rrbracket^{\mathcal{A}}$  and  $\mu(\alpha(\vec{x}_i); r_i^{\mathcal{A}}) > 0$  for all  $i: 1 \leq i \leq k$  is defined as

$$comp(\alpha, G, \mathcal{A}) := \{ comp(\alpha, r_i(\vec{x}_i), \mathcal{A}) | 1 \le i \le k \}$$

where  $\{\ldots\}$  is a multiset.

• The V-complexity of  $\alpha$  for goal  $G = r_1(\vec{x}_1) \& \ldots \& r_k(\vec{x}_k) \& \phi$  in  $\mathcal{A}$  s.t.  $\alpha \in \llbracket \phi \rrbracket_V^{\mathcal{A}}$  and  $\mu(\alpha(\vec{x}_i); r_i^{\mathcal{A}}) > 0$  for all  $i: 1 \leq i \leq k$  is defined as

$$comp_{\mathsf{V}}(\alpha, G, \mathcal{A}) := \min\{comp(\beta, G, \mathcal{A}) \mid \beta \in \llbracket \phi \rrbracket^{\mathcal{A}}, \ \mu(\beta(\vec{x}_i); r_i^{\mathcal{A}}) > 0$$
  
for all  $i : 1 \le i \le k$  and  $\alpha = \beta|_{\mathsf{V}}\}.$ 

The minimum is taken with respect to a total ordering on multisets s.t.  $M \leq M'$  iff  $\forall x \in M \setminus M', \exists x' \in M' \setminus M$  s.t. x < x'.

The following proofs show that the quantitative proof procedure is sound and complete with respect to the above stated semantic concepts. Again, there is a close similarity to the corresponding statements for the non-quantitative case of Höhfeld and Smolka (1988).

**Theorem 3.4 (Soundness).** For each quantitative definite clause specification  $\mathcal{P}_F$ , for each goal G, for each  $\mathcal{L}$  -constraint  $\varphi$ : If there is a quantitative proof tree for G from  $\mathcal{P}_F$  with answer constraint  $\varphi$  and root value v, then  $\varphi_v \to G$  is a logical consequence of  $\mathcal{P}_F$ .

*Proof.* The result is proven by induction on the depth d of the quantitative proof tree, where one unit of depth is from max-node to max-node.

Base: We know that quantitative proof trees of depth d = 0 have to take the form of a single max-node labeled by a satisfiable  $\mathcal{L}$ -constraint  $\psi$  with root value 1. Then  $\psi_1 \rightarrow \psi$  is a logical consequence of  $\mathcal{P}_F$ .

Hypothesis: Suppose the result holds for quantitative proof trees of depth d < n.

Step: Let  $G_0 = r(\vec{x}) \& \phi$  be a goal labeling a quantitative proof tree of depth d = n with answer constraint  $\psi$  and root value h,

let  $G'_0 = q_1(\vec{x}_1) \& \ldots \& q_k(\vec{x}_k) \& \phi \& \phi'$  be a goal labeling the min-node obtained from  $G_0$  via  $\xrightarrow{r}$  using the variant  $C' = r(\vec{x}) \leftarrow_f \phi' \& q_1(\vec{x}_1) \& \ldots \& q_k(\vec{x}_k)$  of a clause C in  $\mathcal{P}_F$ ,

and let  $G_1 = q_1(\vec{x}_1) \& \phi'', \ldots, G_k = q_k(\vec{x}_k) \& \phi''$  be goals labeling max-nodes obtained from  $G'_0$  via  $\stackrel{c}{\longrightarrow}$ .

Then each goal  $G_1, \ldots, G_k$  labels a quantitative proof tree of depth d < n with respective answer constraint  $\psi_1, \ldots, \psi_k$  and root value  $g_1, \ldots, g_k$  s.t.  $h = f \times \min\{g_1, \ldots, g_k\}$  and for each model  $\mathcal{A}$  of  $\mathcal{P}_F : \llbracket \psi \rrbracket^{\mathcal{A}} = \llbracket \psi_1 \& \ldots \& \psi_k \rrbracket^{\mathcal{A}}$ , by definition of min/max tree

- $\implies \psi_1 \ _{g_1} \to G_1, \ldots, \psi_k \ _{g_k} \to G_k$  are logical consequences of  $\mathcal{P}_F$ , by the hypothesis
- ⇒ for each model  $\mathcal{A}$  of  $\mathcal{P}_F$ , for each  $\alpha \in \mathsf{ASS} : \llbracket \psi \rrbracket^{\mathcal{A}} \subseteq \llbracket \phi'' \rrbracket^{\mathcal{A}}$  and if  $\alpha \in \llbracket \psi \rrbracket^{\mathcal{A}}$ , then  $\mu(\alpha(\vec{x}_1); q_1^{\mathcal{A}}) \ge g_1, \ldots, \mu(\alpha(\vec{x}_k); q_k^{\mathcal{A}}) \ge g_k$ , by definition of logical consequence

- $\implies \text{for each model } \mathcal{A} \text{ of } \mathcal{P}_F \text{, for each } \alpha \in \mathsf{ASS} : \llbracket \psi \rrbracket^{\mathcal{A}} \subseteq \llbracket \phi' \rrbracket^{\mathcal{A}} \text{ and if } \alpha \in \llbracket \psi \rrbracket^{\mathcal{A}}, \text{ then} \\ \mu(\alpha(\vec{x}); r^{\mathcal{A}}) \geq f \times \min\{\mu(\alpha(\vec{x}_1); q_1^{\mathcal{A}}), \dots, \mu(\alpha(\vec{x}_k); q_k^{\mathcal{A}})\}, \text{ since each model } \mathcal{A} \text{ of} \\ \mathcal{P}_F \text{ is a model of } C' \text{ iff } \mathcal{A} \text{ is a model of } C$
- $\implies$  for each model  $\mathcal{A}$  of  $\mathcal{P}_F$ , for each  $\alpha \in \mathsf{ASS} : \llbracket \psi \rrbracket^{\mathcal{A}} \subseteq \llbracket \phi \rrbracket^{\mathcal{A}}$  and if  $\alpha \in \llbracket \psi \rrbracket^{\mathcal{A}}$ , then  $\mu(\alpha(\vec{x}); r^{\mathcal{A}}) \geq h$
- $\implies \psi_h \to r(\vec{x}) \& \phi$  is a logical consequence of  $\mathcal{P}_F$ .

The result follows by arithmetic induction.

**Theorem 3.5 (Completeness).** Let  $\mathcal{P}_F$  be a quantitative definite clause specification in  $\mathcal{R}(\mathcal{L})$ ,  $\mathcal{L}$  be closed under renaming,  $\mathcal{A}$  be a minimal model of  $\mathcal{P}_F$ , G be a goal of the form  $r(\vec{x}) \& \phi, \alpha \in \llbracket \phi \rrbracket_V^{\mathcal{A}}$  and  $\mu(\beta(\vec{x}); r^{\mathcal{A}}) = v$  s.t. v > 0 and  $\alpha = \beta |_V$ . Then there exists a quantitative proof tree for G from  $\mathcal{P}_F$  with answer constraint  $\varphi$  and root value v and  $\alpha \in \llbracket \varphi \rrbracket_V^{\mathcal{A}}$ .

*Proof.* The result is proven by induction on  $c = comp_V(\alpha, G, \mathcal{A})$ .

Base: We know that goals with complexity  $c = \emptyset$  have to take the form of a satisfiable  $\mathcal{L}$ constraint  $\chi$ . Then there exists a quantitative proof tree for  $\chi$  from  $\mathcal{P}_F$  consisting of a
single max-node labeled with  $\chi$  and root value 1.

Hypothesis: Suppose the result holds for goals with complexity c < N.

Step: Let  $G_0 = q(\vec{x}) \& \psi, \alpha' \in \llbracket \psi \rrbracket_V^{\mathcal{A}}, \alpha'' \in \llbracket \psi \rrbracket^{\mathcal{A}}, \alpha' = \alpha'' |_V, comp_V(\alpha', G_0, \mathcal{A}) = comp(\alpha'', G_0, \mathcal{A}) = N, comp(\alpha'', q(\vec{x}), \mathcal{A}) := i, \mu(\alpha''(\vec{x}); q^{\mathcal{A}}) = h \text{ and } h > 0.$ 

First we observe, that  $\mu(\alpha''(\vec{x}); q^{\mathcal{A}_i}) = h$ , since  $comp(\alpha'', q(\vec{x}), \mathcal{A}) := i$ 

- $\implies \text{there exists a variant } q(\vec{x}) \leftarrow_f \psi' \& q_1(\vec{x}_1) \& \dots \& q_k(\vec{x}_k) \text{ s.t. } h = f \times \min\{\mu(\alpha(\vec{x}_1); q_1^{\mathcal{A}_{i\perp 1}}), \dots, \mu(\alpha(\vec{x}_k); q_k^{\mathcal{A}_{i\perp 1}})\} \text{ and } \alpha'' \in \llbracket\psi'\rrbracket^{\mathcal{A}_{i\perp 1}} \text{ and } (\mathsf{V} \cup \mathsf{V}(\psi)) \cap \mathsf{V}(\psi' \& q_1(\vec{x}_1) \& \dots \& q_k(\vec{x}_k)) \subseteq \mathsf{V}(q(\vec{x})), \text{ by Definition 3.4} \text{ and renaming closure of } \mathcal{R}(\mathcal{L}), \text{ finite V and infinitely many variables in VAR} \implies G_0 \xrightarrow{r,c} G'_0 \text{ s.t. } G'_0 = q_1(\vec{x}_1) \& \dots \& q_k(\vec{x}_k) \& \psi'' \text{ and } \llbracket\psi''\rrbracket^{\mathcal{A}}_{\mathsf{V}} = \llbracket\psi \& \psi'\rrbracket^{\mathcal{A}}_{\mathsf{V}}, \text{ by definition of the inference rules.}$
- Next,  $\alpha' \in \llbracket \psi'' \rrbracket_{\mathsf{V}}^{\mathcal{A}}$ , since  $\alpha'' \in \llbracket \psi \rrbracket^{\mathcal{A}}$ ,  $\alpha'' \in \llbracket \psi' \rrbracket^{\mathcal{A}_{i\perp 1}} \subseteq \llbracket \psi' \rrbracket^{\mathcal{A}}$ ,  $\alpha'' \in \llbracket \psi \And \psi' \rrbracket^{\mathcal{A}}$ ,  $\llbracket \psi \And \psi' \rrbracket_{\mathsf{V}}^{\mathcal{A}} = \llbracket \psi'' \rrbracket_{\mathsf{V}}^{\mathcal{A}}$  and  $\alpha' = \alpha'' |_{\mathsf{V}}$ .
- Finally,  $comp_{\mathsf{V}}(\alpha', G'_0, \mathcal{A}) < N$ , since  $comp_{\mathsf{V}}(\alpha', G'_0, \mathcal{A}) \leq comp(\alpha'', G'_0, \mathcal{A}) < \{i\} = \{comp(\alpha'', q(\vec{x}), \mathcal{A})\} = comp(\alpha'', G_0, \mathcal{A}) = comp_{\mathsf{V}}(\alpha', G_0, \mathcal{A}) = N.$
- Now we can obtain goals  $G_1 = q_1(\vec{x}_1) \& \psi'', \dots, G_k = q_k(\vec{x}_k) \& \psi''$  from  $G'_0$  s.t.  $\alpha' \in [\![\psi'']\!]_V^{\mathcal{A}}, \, \mu(\alpha''(\vec{x}_1); q_1^{\mathcal{A}}) = g_1 > 0, \dots, \, \mu(\alpha''(\vec{x}_k); q_k^{\mathcal{A}}) = g_k > 0, \, \alpha' = \alpha''|_V$  and  $comp_V(\alpha', G_1, \mathcal{A}) < N, \dots, comp_V(\alpha', G_k, \mathcal{A}) < N$

- $\implies \text{for each goal } G_1, \ldots, G_k, \text{ there exists a quantitative proof tree from } \mathcal{P}_F \text{ with} \\ \text{respective answer constraint } \chi_1, \ldots, \chi_k \text{ and respective root value} \\ g'_1 = g_1, \ldots, g'_k = g_k \text{ and } \alpha' \in [\![\chi_1 \& \ldots \& \chi_k]\!]_V^{\mathcal{A}} = [\![\chi]\!]_V^{\mathcal{A}}, \text{ by the hypothesis} \end{cases}$
- $\implies$  there exists a quantitative proof tree for  $G_0$  from  $\mathcal{P}_F$  with answer constraint  $\chi$  and root value  $h' = f \times \min\{g'_1, \ldots, g'_k\} = f \times \min\{g_1, \ldots, g_k\} = h$  and  $\alpha' \in [\![\chi]\!]_V^{\mathcal{A}}$ .

The result follows by arithmetic induction.

Returning to our toy example, the proof procedure for quantitative definite clause specifications can be illustrated as follows. A min/max derivation tree for the query q(X) & X = e and the program of Fig. 3.1 is given in Fig. 3.3.

$$q(X) \& X = e$$
  

$$\max\{.7\} = .7$$
  

$$|r$$
  

$$1, p(X) \& X = e$$
  

$$1 \times \min\{.7\} = .7$$
  

$$|c$$
  

$$p(X) \& X = e$$
  

$$\max\{.7,.5\} = .7$$
  

$$r$$
  

$$2, X = e \& X = a$$
  

$$.7 \times \min\{1\}$$
  

$$|c$$
  

$$X = a$$
  

$$1$$
  

$$X = b$$
  

$$1$$
  

$$1$$
  

$$X = e$$
  

$$X = b$$
  

$$1$$
  

$$x = a$$
  

$$x = b$$
  

$$1$$
  

$$x = a$$
  

$$x = b$$
  

$$1$$
  

$$x = a$$
  

$$x = b$$
  

$$x =$$

Figure 3.3: Min/max tree for quantitative constraint logic program

This tree contains two success nodes, X = a and X = b, from which two distinct quantitative proof trees can be obtained (see Fig. 3.4).

Soundness of quantitative CLP tells us that corresponding to the quantitative proof tree with answer constraint X = a (resp. X = b) and root value .7 (resp. .5), we know that the quantitative formula  $X = a_{.7} \rightarrow q(X) \& X = e$  (resp.  $X = b_{.5} \rightarrow q(X) \& X = e$ ) is a logical consequence of the program of Fig. 3.1. This can easily be verified from the minimal model constructed in Fig. 3.2.

Completeness says that for an object  $\langle \alpha(X) \rangle$  assigned by  $\alpha \in [\![X = e]\!]^{\mathcal{I}}$  with membership degree  $\mu(\langle \alpha(X) \rangle; \mathbf{q}^{\mathcal{A}}) = .7$  to the denotation of  $\mathbf{q}$  under the minimal model  $\mathcal{A}$ , we have a corresponding proof tree with answer constraint X = a and root value .7 and  $\alpha \in [\![X = a]\!]^{\mathcal{I}}$ .

$$q(X) \& X = e$$
 $q(X) \& X = e$ 
 $\max\{.7\}=.7$ 
 $\max\{.5\}=.5$ 
 $|r$ 
 $|r$ 
 $1, p(X) \& X = e$ 
 $1, p(X) \& X = e$ 
 $1 \times \min\{.7\}=.7$ 
 $1 \times \min\{.5\}=.5$ 
 $|c$ 
 $p(X) \& X = e$ 
 $\max\{.7\}=.7$ 
 $p(X) \& X = e$ 
 $\max\{.7\}=.7$ 
 $p(X) \& X = e$ 
 $\max\{.7\}=.7$ 
 $p(X) \& X = e$ 
 $\max\{.5\}=.5$ 
 $|r$ 
 $|r$ 
 $2, X = e \& X = a$ 
 $.7 \times \min\{1\}$ 
 $.5 \times \min\{1\}$ 
 $|c$ 
 $|c$ 
 $X = a$ 
 $X = b$ 
 $1$ 
 $1$ 

Figure 3.4: Quantitative proof trees for quantitative constraint logic program

Similarly, for an object  $\langle \alpha'(X) \rangle$  with  $\alpha' \in [\![X = e]\!]^{\mathcal{I}}$  and  $\mu(\langle \alpha'(X) \rangle; q^{\mathcal{A}}) = .5$ , we have a proof tree with answer constraint X = b and root value .5 and  $\alpha' \in [\![X = b]\!]^{\mathcal{I}}$ .

### 3.6 Parsing and Searching in Quantitative CLGs

The quantitative CLP scheme presented in the last chapter allows for a definition of the parsing problem (and similarly of the generation problem) for quantitative CLGs in the following way: Given a program  $\mathcal{P}_F$  (encoding some quantitative CLG) and a query G (encoding some input sentence), we ask if we can infer a  $\mathcal{P}_F$  -answer  $\varphi$  of G (encoding a parse of the input sentence) at a value v (encoding the weight of the parse) proving  $\varphi_v \to G$  to be a logical consequence of  $\mathcal{P}_F$ . That is, according to the soundness and completeness results presented above, the operational concept of a quantitative proof tree has a declarative counterpart in the form of a quantitative  $\mathcal{P}_F$  -answer. Truth-functionally, a quantitative  $\mathcal{P}_F$  -answer tells us that the answer constraint  $\varphi$  contributes a truth-value v to the goal G in every model of  $\mathcal{P}_F$ . In terms of membership values, this means that a  $\mathcal{P}_F$  -answer to a query  $G = r(\vec{x}) \& \phi$  at value v is a satisfiable  $\mathcal{L}$  -constraint  $\varphi$  such that for each model  $\mathcal{A}$  of  $\mathcal{P}_F$  holds: If  $\varphi$  is satisfiable, then  $\phi$  is satisfiable and all objects assigned to  $\vec{x}$  by a solution of  $\varphi$  are in the denotation of  $r(\vec{x})$  at a membership value of at least v.

#### 3.6.1 Quantitative Feature-Based CLGs

Returning to the simple linguistic CLG of Fig. 2.7, the formal scheme described above can be illustrated as follows.

The quantitative constraint logic program  $\mathcal{P}_F$  of Fig. 3.5 is obtained from the program of Fig. 2.7 simply by adding numerical factors to the program clauses.

1 phrase(X)  $\leftarrow_{f_1} X = (phrase \wedge CAT : s \wedge DTR1:CAT : n \wedge DTR2:CAT : v \wedge DTR1:AGR : Y \wedge DTR2:AGR : Y \wedge DTR1 : Z_1 \wedge DTR2 : Z_2) \& \operatorname{sign}(Z_1) \& \operatorname{sign}(Z_2).$ 

2 phrase(X)  $\leftarrow_{f_2} X = (phrase \land CAT : np \land DTR1:CAT : n \land DTR2:CAT : n \land DTR1 : Z_1 \land DTR2 : Z_2) \& \operatorname{sign}(Z_1) \& \operatorname{sign}(Z_2).$ 

 $3 \operatorname{word}(X) \leftarrow_{f_3} X = (word \wedge \operatorname{CAT} : n \wedge \operatorname{PHON} : Clinton \wedge \operatorname{AGR} : sg).$ 

4 word(X)  $\leftarrow_{f_4} X = (word \land CAT : v \land PHON : talks \land AGR : sg).$ 

5 word(X)  $\leftarrow_{f_5} X = (word \wedge CAT : n \wedge PHON : talks \wedge AGR : pl).$ 

6 sign $(X) \leftarrow_{f_6}$  phrase(X).

7  $\operatorname{sign}(X) \leftarrow_{f_7} \operatorname{word}(X)$ .

Figure 3.5: Quantitative feature-based constraint logic grammar

Given the quantitative CLG of Fig. 3.5 and a goal G of the form

 $X = (sign \land DTR1: PHON : Clinton \land DTR2: PHON : talks) \& sign(X),$ 

again encoding the input sentence *Clinton talks*, we can infer two different proof trees for G, each with a specific answer constraint, encoding a parse, and a specific root value, encoding the preference value of the parse. Again, we will depict only success branches and consider the constraint solver as a black box. The two derivations are shown in Figs. 3.6 and 3.7.

The answer constraint  $\phi$  of the first derivation is obtained by constraint solving of the terminal constraints of the first proof tree. We get

 $\star [\dots] \& \star [\dots] \xrightarrow{c} X = (phrase \land CAT : s \land DTR1 : word \land DTR1 : CAT : n \land DTR1 : PHON : Clinton \land DTR1 : AGR : Y \land DTR1 : AGR : sg \land DTR2 : word \land DTR2 : CAT : v \land DTR2 : PHON : talks \land DTR2 : AGR : Y \land DTR2 : AGR : sg)$ 

yielding the reading  $[Clinton_N \ talks_V]_S$  with weight

$$\upsilon = f_6 \times f_1 \times min\{f_7 \times f_3, f_7 \times f_4\}.$$

The answer constraint  $\psi$  of the second derivation is



Figure 3.6: Quantitative derivation of  $[Clinton_N \ talks_V]_S$ 

 $X = (sign \land DTR1: PHON : Clinton \land DTR2: PHON : talks) \& sign(X)$  $\mathbf{f_6}\times\mathbf{f_2}\times\mathbf{min}\{\mathbf{f_7}\times\mathbf{f_3},\mathbf{f_7}\times\mathbf{f_5}\}$ r6,  $X = (sign \land DTR1: PHON : Clinton \land DTR2: PHON : talks) \& phrase(X)$  $\mathbf{f_6} \times \mathbf{f_2} \times \min\{\mathbf{f_7} \times \mathbf{f_3}, \mathbf{f_7} \times \mathbf{f_5}\}$ c $X = (sign \land DTR1: PHON : Clinton \land DTR2: PHON : talks) \& phrase(X)$  $\mathbf{f_2} \times \min\{\mathbf{f_7} \times \mathbf{f_3}, \mathbf{f_7} \times \mathbf{f_5}\}$ |r|2,  $X = (sign \land DTR1: PHON : Clinton \land DTR2: PHON : talks)$ &  $X = (phrase \land CAT : np \land DTR1: CAT : n \land DTR2: CAT : n$  $\wedge$  DTR1 :  $Z_1 \wedge$  DTR2 :  $Z_2$ )  $\& \operatorname{sign}(Z_1) \& \operatorname{sign}(Z_2)$  $\mathbf{f_2} \times \min\{\mathbf{f_7} \times \mathbf{f_3}, \mathbf{f_7} \times \mathbf{f_5}\}$ cc $*[X = (phrase \land CAT : np \land DTR1 : word \land DTR1 : CAT : n$  $*[...] \& sign(Z_2)$  $\land$  DTR1: PHON : Clinton  $\land$  DTR2 : word  $\land$  DTR2: CAT : n  $f_7 \times f_5$  $\wedge$  DTR2: PHON : talks  $\wedge$  DTR1 :  $Z_1 \wedge$  DTR2 :  $Z_2$ )] |r| $\& \operatorname{sign}(Z_1)$ 7, \* [...] & word( $Z_2$ )  $f_7\,\times f_3$  $f_7 \times f_5$ r|c|7,  $*[X = (phrase \land CAT : np \land DTR1 : word \land DTR1 : CAT : n]$  $*[\dots]$  & word $(Z_2)$  $\land$  DTR1: PHON : Clinton  $\land$  DTR2 : word  $\land$  DTR2: CAT : n  $f_5$  $\land$  DTR2: PHON : talks  $\land$  DTR1 :  $Z_1 \land$  DTR2 :  $Z_2$ )] r& word( $Z_1$ ) 5, \*[...]  $f_7 \times f_3$  $\& Z_2 = (word \land CAT : n)$ c $\wedge$  PHON : talks  $\wedge$  AGR : pl)  $*[X = (phrase \land CAT : np \land DTR1 : word \land DTR1 : CAT : n$  $\mathbf{f}_{5}$  $\land$  DTR1: PHON : Clinton  $\land$  DTR2 : word  $\land$  DTR2: CAT : n |c| $\wedge$  DTR2: PHON : talks  $\wedge$  DTR1 :  $Z_1 \wedge$  DTR2 :  $Z_2$ )]  $\dagger [X = (phrase \land CAT : np$  $\& \operatorname{word}(Z_1)$  $\land$  DTR1 : word  $\land$  DTR1: CAT : n  $f_3$  $\land$  DTR1: PHON : Clinton r $\land$  DTR2 : word  $\land$  DTR2: CAT : n 3,  $*[X = (phrase \land CAT : np \land DTR1 : word \land DTR1 : CAT : n]$  $\wedge$  DTR2: PHON : talks  $\land$  DTR1: PHON : Clinton  $\land$  DTR2 : word  $\land$  DTR2: CAT : n  $\wedge$  DTR2: AGR : pl]  $\wedge$  DTR2: PHON : talks  $\wedge$  DTR1 :  $Z_1 \wedge$  DTR2 :  $Z_2$ )] 1 &  $Z_1 = (word \land CAT : n \land PHON : Clinton)$  $f_3$ c $\star[X = (phrase \land \text{ CAT} : np \land \text{ DTR1} : word \land \text{ DTR1} : \text{ CAT} : n$  $\land$  DTR1: PHON : Clinton  $\land$  DTR1: AGR : sg  $\land$  DTR2 : word  $\wedge$  DTR2: CAT :  $n \wedge$  DTR2: PHON : talks)] 1

Figure 3.7: Quantitative derivation of  $[Clinton_N talks_N]_{NP}$ 

 $\star [\dots] \& \dagger [\dots] \xrightarrow{c} X = (phrase \land CAT : np \land DTR1 : word \land DTR1 : CAT : n \land DTR1 : PHON : Clinton \land DTR1 : AGR : sg \land DTR2 : word \land DTR2 : CAT : n \land DTR2 : PHON : talks \land DTR2 : AGR : pl)$ 

yielding the reading  $[Clinton_N \ talks_N]_{NP}$  with weight

 $\tau = f_6 \times f_2 \times min\{f_7 \times f_3, f_7 \times f_5\}.$ 

Suppose now that we have a subjective weight assignment for the factors of the quantitative CLG of Fig. 3.5 where  $f_1 > f_2$  and  $f_4 > f_5$ . That is, we prefer the rule  $S \to N V$  over the rule  $NP \to N N$  to describe a phrase. Furthermore, the terminal rule  $V \to talks$ , encoding the word talks as a verb, is preferred over the rule  $N \to talks$ , encoding it as a noun. Clearly, we get a preference of the answer constraint  $\phi$ , encoding the reading  $[Clinton_N talks_V]_S$ , over the answer constraint  $\psi$ , encoding the reading  $[Clinton_N talks_N]_{NP}$ , with  $v > \tau$ .

#### 3.6.2 Alpha-Beta Searching in Quantitative CLGs

As proposed by van Emden (1986), search strategies such as alpha-beta pruning that are well known in game theory can be used quite directly to define efficient search strategies for quantitative rule sets. The same technique can be applied to the proof procedure of quantitative CLP. Alpha-beta pruning is a technique to speed up the search in min/max trees without loss of information. For our application, alpha-beta pruning can be used to efficiently search a min/max derivation tree for the maximal valued proof tree. The fact that no information is lost in alpha-beta pruning means in our context that the maximal valued proof tree is guaranteed to be found. Furthermore, in general, the amount of search needed to find the best proof for a goal, i.e. the maximal valued proof tree for a goal from a program, will be reduced remarkably by controlling the search by the alpha-beta algorithm.

The central concepts of alpha-beta pruning can be summarized as follows (see Nilsson (1982)).

Usually some form of depth-first search is employed in alpha-beta pruning. The search procedure associates with each max-node (resp. min-node) a dynamic alpha-value (resp. beta-value). These values are based on the static values of terminal nodes and will be backed-up in subsequent search by lookahead in the tree.

The search procedure starts with a maximum depth execution of depth-first search, initializing the alpha and beta values of the first subtree. During search, alpha and beta values are computed as follows:

- The alpha value of a max-node is the maximum of the current values of its successors.
- The **beta value** of a min-node is the minimum of the current values of its successors, multiplied by the factor of the clause labeling the min-node.

The rules for discontinuing the search are as follows:

- Alpha-cutoff: Search can be discontinued below any min-node having a beta value less than or equal to the alpha value of any of its max-node ancestors. The final backed-up value of this min-node can then be set to its beta value.
- **Beta-cutoff**: Search can be discontinued below any max-node with the product of its alpha value and the factor of the rule labeling its min-node ancestor being greater than or equal to the beta value of this min-node ancestor for all min-node ancestors. The final backed-up value of this max-node can then be set to its alpha value.

The procedure terminates when all of the successors of the root node have been given a final backed-up value. The maximal valued proof tree is then the one taking as single successor of each of its max-nodes the successor with the maximal final backed-up value. This proof tree is found efficiently if the original min/max tree can be pruned by the alpha-beta procedure to a tree consisting of a relatively small number of nodes.

Let us illustrate these concepts with a simple example. A sample artificial program is given in Fig. 3.8.

1 
$$p(X) \leftarrow_{.7} r(X) \& s(X).$$
  
2  $r(X) \leftarrow_{.8} X = a.$   
3  $s(X) \leftarrow_{.9} X = a.$   
4  $s(X) \leftarrow_{.2} r(X).$   
5  $p(X) \leftarrow_{.7} t(X) \& r(X) \& s(X)$   
6  $t(X) \leftarrow_{.1} X = a.$ 

Figure 3.8: Quantitative constraint logic program

The complete min/max derivation tree for the query p(X) & X = a to the program of Fig. 3.8 is given in Fig. 3.9.

The concept of alpha-beta pruning can be illustrated with this example as follows (see Fig. 3.10). The alpha value  $\alpha = .9$  of the max-node  $\mathbf{s}(X) \& X = a$  times the factor .7 of the minnode ancestor is greater than the beta value  $\beta = .56$  of this min-node. Since we know that this alpha value cannot be decreased by further evaluation of the subtrees of this max-node, and since we are interested in the minimum of the values of the successors of this min-node, we can cut off the search below this max-node without a risk of losing information relevant to the final maximal valued proof tree. This cutoff is indicated by the dotted line below this max-node in Fig. 3.10. In a similar way, search below the min-node 5,  $\mathbf{t}(X) \& \mathbf{r}(X) \& \mathbf{s}(X) \& X = a$  can be discontinued because the non-decreasing beta value  $\beta = .07$  of this node is already smaller than the alpha value  $\alpha = .56$  of its max-node ancestor. The pruning of the two subtrees of this



Figure 3.9: Complete search of a quantitative derivation tree

min-node again is indicated by dotted lines in Fig. 3.10. Again, there is no risk of information loss in this pruning step.

Clearly, in each application of the alpha-beta procedure, the number of nodes to be generated and evaluated is minimal when the number of cutoffs is maximal. The best case occurs when the maximal valued proof tree is reached first in the depth-first search. In the worst case, no gain in search efficiency is obtained at all, i.e., all nodes of the min/max tree have to be generated. In either case, the maximal valued proof tree is guaranteed to be left unpruned. Risking loss of relevant information, the alpha-beta procedure can be improved by setting an initial alpha value for the root note which allows to cut off search branches with root value lower than this initial value. For a thorough analysis of the properties of alpha-beta pruning the reader is referred to Knuth and Moore (1975).

Furthermore, it should be noted that a strict application of alpha-beta pruning is possible only for quantitative CLP based on min/max trees. Suppose for example that the minimum operator is replaced by a product operator throughout the declarative as well as operational semantics of quantitative CLP. This replacement could be motivated by the aim to consider the contribution of all instead only one antecedent atom to the weight of the consequent. To efficiently search for the maximal valued proof tree in such a setting, a version of alpha-beta pruning employing only alpha-cutoffs has to be used. In this setting, additional beta-cutoffs can improve the search efficiency for finding a *good* proof tree, but possibly cut off parts of the *best* proof tree, i.e., here attention has to be paid to the risk of losing information relevant



Figure 3.10: Alpha-beta search of quantitative derivation tree

to the maximal valued proof tree.

# 3.7 Summary and Discussion

In this chapter we presented a formal framework for quantitative CLP. In this framework CLP clauses were attached by arbitrary numerical weights. Such weighted clauses were interpreted in a model-theory based on concepts of fuzzy set algebra. The quantitative system was shown to be sound and complete with respect to a fixpoint semantics based on minimal models in this model-theory. We illustrated the concepts of quantitative CLP by a simple quantitative feature-based CLG. Furthermore, we showed how to adapt the search algorithm of alpha-beta pruning to searching efficiently for the highest weighted proof tree in a quantitative CLP system.

The advantage of quantitative CLP clearly is the freedom it offers to the grammar writer or implementer to specify arbitrary weights in a formally clear and efficient programming framework. Such weights could be specified, e.g., as subjective preference values (Erbach 1993b, Kim 1994, Douglas and Dale 1992), subjective values expressing graded grammaticality (Erbach 1993a), or subjective probabilities (Erbach 1998). Calculation with arbitrary such weights can be interpreted in a unique well-defined formal framework. Furthermore, generalizations of the formal system to particular applications which require particular calculation schemes are easily obtainable. For example, if we want to model probabilistic context-free grammars (Booth and Thompson 1973) in quantatitive CLP, we simply must attach subjective probability measures to a context-free program according to the conditions of Booth and Thompson (1973), and replace the minimum operator by a product operator in all relevant definitions of the declarative and operational semantics of quantitative CLP. Unfortunately, such changes in the weight calculation model prevent a direct application of alpha-beta pruning for efficient disambiguation. For example, in the case of probabilistic context-free CLP, only a restricted version of alpha-beta pruning using exclusively alpha-cutoffs is applicable. Alternatively, one could use a form of best-first pruning for the search task, i.e., a search rule which selects the highest weighted clause at each derivation step. Clearly, this approach does not guarantee that the highest weighted proof tree is found, but offers only an approximate heuristic search procedure.

However, regardless of the specific choice of weights, a proper specification of a multitude of weights can be very complex and is always user-dependent. In several applications, one would like to trade in the flexibility of subjective weight assignment for automatic and reusable methods for estimating weights from empirical data. One solution to this problem is to use automatic methods for statistical inference to induce values of probabilistic parameters from empirical data.

In the next chapter, we present an framework of probabilistic CLP which addresses the problem of finding a proper probability distribution over the set of proof trees of a constraint logic program and of using statistical estimation methods to infer parameters from empirical data. Clearly, even if it would be possible to specify a model-theoretic semantics for such a system, it is superfluous to do so in the context of automatic statistical inference. Rather, the interest is here in the stochastic semantics of CLP provided by the probabilistic and statistical methods used.
# Chapter 4

# Probabilistic CLP: Probabilistic Modeling and Statistical Inference from Incomplete Data

In this chapter we present a probabilistic model for CLP and a novel method for statistical inference of the parameter values of such a model from incomplete training data. We show mononoticity and convergence of the new algorithm to the desired maximum likelihood estimates. Furthermore, we show the usefulness of the statistical approach by a small-scale experiment on estimating feature-based CLGs. We present a novel algorithm to infer the properties of such parametric probability models from incomplete data and discuss different approaches for approximate computation for the inference task. Moreover, we discuss the possibilities of using the structure of the probabilistic model to guide the search in finding the most probable proof tree in probabilistic CLP and present as heuristic search method for this task.

This chapter is based upon work previously published in Riezler (1997), Riezler (1998a), Riezler (1998b), and Johnson, Geman, Canon, Chi, and Riezler (1999).

## 4.1 Introduction and Overview

In the previous chapter we presented a formal semantics for a system of quantitative CLP. This formal semantics and the connected quantitative inference system were crucially based upon open parameters for subjective weights. Most approaches to probabilistic logic programming interpret such weights as subjective probabilities, and concentrate on inference systems and formal semantics for programming systems with user-defined probabilities attached to the formulae of the language. The aim of such approaches is the development of sound and complete logic programming systems where the handling of weights is restricted to accord to the laws of probability theory. That is, these approaches aim to connect logical inference with probabilistic inference.

In this chapter, we present a completely different approach to probabilistic CLP. In this approach, subjective assignment of probabilities is replaced by automatic and reusable methods for estimating empirical probabilities from data. The central aims of this approach are the specification of a probability distribution over the proof trees given by a program, and the use of statistical methods to infer the values of the probabilistic parameters from empirical data. That is, in this setting, the weight of a CLP proof tree is determined directly by a probability distribution over proof trees rather than by quantitative calculation scheme referring to weighted clauses. The parameters of the probability distribution are determined by statistical inference from empirical data rather than by an assignment of subjective weights to clauses. Furthermore, the specific properties of the parametric probability model can be inferred by statistical methods. That means, in this chapter we do not only turn from quantitative to probabilistic inference but, what is more, to statistical inference. In such a setting, the connection of probability theory, semantic fixpoint theory and logical inference theory is not of interest since the specification of probabilistic parameters is done by automatic statistical methods and not manipulable by the user. Rather, we are interested in the stochastic semantics defined by the methods of probabilistic modeling and statistical inference.

The statistical problem we consider here is the problem of statistical parameter estimation. We assume that the statistical properties of a given sample of observations  $O = O_1, \ldots, O_n$ can be described by a parametric family of probability distributions. That is, the probability distribution that generated the data is assumed to be completely known except for the values of a vector  $\theta$  of parameters. We then ask how the unknown value of  $\theta$  can be estimated from the observation sequence O, i.e., a statistical inference is made about the values of the parameters defining that family. Recent interest in statistical approaches to NLP can be attributed to the fact that solutions to such statistical problems can lead quite directly to effective, but conceptually simple and mathematically clear solutions to various problems in NLP. In the context of structural ambiguity resolution in NLP systems, this connection is as follows: Given a probabilistic grammar depending on parameter vector  $\theta$  and given a training corpus O, a solution  $\hat{\theta}$  to the parameter estimation problem will adapt the model parameters to best account for the input corpus. This tuning of the grammar to a particular natural language corpus is a necessary prerequisite for probabilistic disambiguation. That is, when the plausibility of a parse is connected with its probability, the assumption that the correct parse of a sentence is its most probable parse can be made with some justification if the underlying probabilistic grammar is based on parameter values  $\hat{\theta}$  estimated from large data sets of natural language.

The aim of this chapter is to solve open problems in statistical inference and probabilistic modeling of constraint-based grammars. Following Abney (1996), we choose the parametric

family of log-linear probability distributions to model such grammars. The great advantage of log-linear models is their generality and flexibility. Log-linear models allow to describe arbitrary context dependencies in the data by choosing a few salient properties of the data as the defining properties of the model. In contrast to most approaches to probabilistic grammars, with log-linear models we are not restricted to build our models on production rules or other configurational properties of the data. Rather, we have the virtue of employing essentially arbitrary properties in our models. For example, heuristics on preferences of grammatical functions or on attachment preferences as used in Srinivas, Doran, and Kulick (1995), or the preferences in lexical relations as used in Alshawi and Carter (1994) can be integrated into a log-linear model very easily. However, the step from simple rule-based probability models to general log-linear models requires also a more general and more complex estimation algorithm. The estimation algorithm for log-linear models used by Abney (1996) is the iterative scaling method of Della Pietra, Della Pietra, and Lafferty (1997). This algorithm allows to recast the optimization of weights of preference functions as done by Srinivas, Doran, and Kulick (1995) or Alshawi and Carter (1994) as estimation of parameters associated with the properties of a log-linear model. However, there is a drawback: In contrast to rule-based models where efficient estimation algorithms from incomplete, i.e., unannotated data exist, the iterative scaling estimation method of Della Pietra, Della Pietra, and Lafferty (1997) applies only to complete, i.e., fully annotated training data. Unfortunately, the need to rely on large samples of complete data is impractical. For parsing applications, complete data means several personvears of hand-annotating large corpora with specialized grammatical analyses. This task is always labor-intensive, error-prone, and restricted to a specific grammar framework, a specific language, and a specific language domain.

Thus, the first open problem to solve is fo find automatic and reusable techniques for parameter estimation of probabilistic constraint-based grammars from incomplete data. We will present a general estimation algorithm for log-linear models from incomplete data which can be seen as an extension of the iterative scaling method of Della Pietra, Della Pietra, and Lafferty (1997). We prove monotonicity and convergence of the new algorithm to (local) maxima of the incomplete-data log-likelihood function, and show how automatic property selection can be done from incomplete data.

A further open problem is the empirical evaluation of the performance of probabilistic constraint-based grammars in terms of finding human-determined correct parses. We present an experiment with a log-linear model employing a few hundred general properties encoding grammatical functions, attachment preferences, branching behaviour, parallelism, and other general properties of constraint-based parses. The experiment was conducted on a small scale but clearly shows the usefulness of general properties in order to get good results in a linguistic evaluation.

Clearly, for larger scales, problems arise concerning the tractability of the estimation for-

mulae. We discuss the applicability of several approximation methods to our problem of statistical inference from incomplete data, including Newton's method, Monte Carlo methods, or methods for approximating expectations via pseudo-likelihood approaches.

A further open problem is the efficient search for most probable parses, i.e., best-parse search, in parsing systems based on probabilistic constraint-based grammars. Instead of listing all possible parses and selecting the most probable one, one would like to use the structure of the probabilistic model to guide the search for the most probable analysis. Most popular approaches use the search technique of the Viterbi algorithm (Viterbi (1967), Forney (1973)) to solve this problem, but there is as yet no solution for probabilistic constraint-based grammars. We show that standard methods for best-parse search are only of limited use for probabilistic models involving context-dependencies, and make the move to approximate heuristic methods.

To summarize, our approach satisfies the following requirements. It

- is generally applicable to probability models involving context-dependencies, and especially to a probabilistic model for CLP over arbitrary constraint languages,
- provides automatic and reusable techniques for statistical inference from incomplete data for such probability models, and
- is accompanied with search techniques for finding most probable analyses in probabilistic CLP.

This chapter is organized as follows. Is Sect. 4.2 we discuss related previous approaches to statistical inference for probabilistic constraint-based grammars.

In Sect. 4.3 we introduce the basic concepts of maximum likelihood estimation from incomplete data via the EM algorithm.

Sect. 4.4 discusses the problem of applying a popular instance of this algorithm, namely Baum's maximization technique for stochastic context-free models, to parameter estimation for probabilistic CLP.

Sect. 4.5 and 4.6 present in detail a solution to this problem by introducing a log-linear probability model for CLP coupled with an incomplete-data inference algorithm for such models. This section includes a detailed proof of monotonicity and convergence of the inference algorithm.

Sect. 4.7 presents an empirical evaluation of the applicability of general log-linear distributions to probabilistic constraint-based grammars in a small-scale experiment on estimating a log-linear model on constraint-based parses.

Sect. 4.8 discusses computation issues such as the use of Monte Carlo methods, Newton's numerical method, and other approximation techniques in the context of this inference process.

Sect. 4.9 discusses the applicability of standard parsing and search methods to contextdependent constraint-based models, and presents a heuristic method for searching for best parses in CLGs.

## 4.2 Previous Work

An approach to define estimators for probabilistic constraint-based grammars which has been applied to nearly all constraint-based formalisms is a renormalized extension of the estimator for stochastic regular (Baum, Petrie, Soules, and Weiss 1970) or context-free grammars (Baker 1979) to constraint-based models. Examples for this approach are, e.g., stochastic unification-based grammars (Briscoe and Waegner 1992; Briscoe and Carroll 1993), stochastic constraint logic programming (Eisele 1994), stochastic head-driven phrase structure grammar (Brew 1995), stochastic logic programming (Miyata 1996), stochastic categorial grammars (Osborne and Briscoe 1997) or data-oriented approaches to lexical-functional grammar (Bod and Kaplan 1998). Since the estimation technique for context-free models is based on the assumption of mutual independence of the model's derivation steps, but context-dependent constraints on derivations are inherent to constraint-based grammars, a loss in probability mass due to failure derivations is caused in these approaches. However, the necessary renormalization of the probability distribution on derivations with respect to consistent derivations causes a general deviance of the resulting estimates from the desired maximum likelihood estimates. This was shown firstly by Abney (1996) for estimation of constraint-based models from complete data. We will make a similar argument for incomplete data in the following. Optimization-theoretically these approaches can be described as maximization procedures for pseudo-likelihood functions for context-free models where the probability distribution on context-free derivations is restricted to consistent derivations in the constraint-based sense. Maximum pseudo-likelihood estimators for context-free models certainly are sensible, e.g., if the aim is to constrain an inherently context-free language to include only linguistically plausible derivations as is done by introducing bracketing constraints on context-free derivations by Pereira and Schabes (1992). However, it is questionable if is the best way to model constraintbased grammars probabilistically by context-free models which respect constraints only indirectly to discard derivations. The move to log-linear models as is done in our approach clearly has several advantages. Since there is linguistically no reason to base probabilistic grammars on rule-properties, we can now exploit the flexibility of log-linear distributions and model the context-dependencies in the data directly. Furthermore, since the new family of parametric probability models requires new estimation techniques, we can again take consistent maximum likelihood estimators as the optimization procedures of our choice.

Other approaches to probabilistic constraint-based models have been presented which define custom-built statistical inference procedures for specialized parsing models including a limited amount of context-dependency. For example, the model presented by Goodman (1998) conditions on a finite set of categorial features beyond the nonterminal of each node which makes it possible to explicitly unfold the dependencies in the parsing model. This allows for the use of standard dynamic programming techniques for computation. In the approaches of Magerman (1994) and Ratnaparkhi (1998) general statistical inference methods, namely decision trees and maximum-entropy methods, are used to infer weights associated to the actions of specialized parsing models including limited context-dependency. However, it is difficult to generalize these models to arbitrary log-linear models on constraint-based grammars, concerning both the choice of properties and the issue of efficient computation. Clearly, a careful choice of properties makes it possible to tune specialized models to maximum accuracy and efficiency, which does not hold for the general case<sup>1</sup> The aim of our approach is to address problems concerning estimation, property design, or approximation methods for general log-linear models and show these general ideas to be applicable in practice.

## 4.3 Maximum Likelihood Estimation from Incomplete Data via the EM Algorithm

A constant companion during the course of this chapter will be the statistical estimation technique of the Expectation-Maximization (EM) algorithm. The fact that both Baum's estimation technique, which is shown not to be applicable to probabilistic CLP in Sect. 4.4.2, and the incomplete-data estimation algorithm for log-linear models we present in Sects. 4.5-4.8, can be seen as instances of the EM algorithm, justifies a closer look at this estimation scheme.

#### 4.3.1 General Theory of the EM Algorithm

The EM algorithm has been introduced by Dempster, Laird, and Rubin (1977), although central parts of the general theory can be found earlier in special applications, e.g., in Baum, Petrie, Soules, and Weiss (1970). Various applications and extensions of the algorithm are discussed in Little and Rubin (1987) and, more recently, in McLachlan and Krishnan (1997).

The EM method is a technique for maximum likelihood estimation (MLE) from incomplete data. For a parametric family of probability distributions depending on parameter vector  $\theta$ and a given sample of training data from this parametric family, MLE defines the estimate  $\hat{\theta}$ of  $\theta$  as a value of  $\theta$  which maximizes the likelihood of the training sample. MLE from observed complete data is particularily easy for many statistical problems, thanks to the nice form of the complete-data (log-)likelihood function. The problem the EM algorithm especially addresses is the case where the observed data are incomplete. That is, we observe only a function of

<sup>&</sup>lt;sup>1</sup>For example, as noted by Goodman (1998), the computational complexity of his dynamic programming algorithm for probabilistic feature-grammars is exponential in the general case.

complete data, which themselves are unobserved. Because of this indirect, hidden character of the complete data, MLE from incomplete data is difficult.

In the following, an incomplete-data estimation setting is assumed to consist of

- a sample space  $\mathcal{Y}$  of observed, incomplete data,
- a sample space  $\mathcal{X}$  of unobserved, complete data,
- a many-to-one function Y : X → Y s.t. Y(x) = y is the unique observation corresponding to the complete datum x, and its inverse X : Y → 2<sup>X</sup> s.t. X(y) = {x|Y(x) = y} is the countably infinite set of complete data corresponding to the observation y,
- a complete-data specification  $p_{\theta}(x)$  with parameters  $\theta \in \Theta$ ,
- an incomplete data specification  $g_{\theta}(y)$  which is related to the complete-data specification by marginalization as

$$g_{\theta}(y) = \sum_{x \in X(y)} p_{\theta}(x).$$

Let  $y_1, y_2, \ldots, y_N$  be a random sample from  $\mathcal{Y}$ , i.e., values of independently and identically distributed (i.i.d.) random variables on  $\mathcal{Y}$ . Let  $p[f] = \sum_{\omega \in \Omega} p(\omega) f(\omega)$  denote the expectation of a function  $f: \Omega \to \mathbb{R}$  with respect to a probability distribution p on  $\Omega$ . If f is a multivariable function  $f(\omega', \omega)$ , then the expectation of f with respect to  $p(\omega)$  is written  $p[f(\omega', \cdot)]$ . Furthermore, let the empirical probability  $\tilde{p}(y)$  of an incomplete data type be defined as  $\tilde{p}: \mathcal{Y} \to \mathbb{R}$  s.t.  $\tilde{p}(y) = N^{\perp 1} \sum_{i=1}^{N} \delta_{y_i,y}$  where the Kronecker delta  $\delta_{y_i,y} = \begin{cases} 1 & \text{if } y_i = y, \\ 0 & \text{otherwise.} \end{cases}$  Then the incomplete-data log-likelihood L is defined for a random sample from  $\mathcal{Y}$  as a function of  $\theta$  as

$$L(\theta) = \ln \prod_{y \in \mathcal{Y}} g_{\theta}(y)^{\tilde{p}(y)} = \sum_{y \in \mathcal{Y}} \tilde{p}(y) \ln g_{\theta}(y) = \tilde{p}[\ln g_{\theta}].$$

The EM algorithm is directed at finding a value  $\hat{\theta}$  of  $\theta \in \Theta$  that maximizes L as a function of  $\theta$  for a given random sample from  $\mathcal{Y}$ , i.e.,

$$\hat{\theta} = \underset{\theta \in \Theta}{\operatorname{arg\,max}} L(\theta) \text{ where } L(\theta) = \tilde{p}[\ln g_{\theta}] = \tilde{p}[\ln \sum_{x \in X(\cdot)} p_{\theta}(x)]$$

The summation inside this logarithm can make MLE from incomplete data difficult even when complete-data MLE is easy.

The old idea formalized in the EM algorithm can be stated informally as follows: 1. Replace unobserved data values by expected values, 2. perform MLE from the expected complete data, 3. recompute the unobserved data expectations using the new parameter estimates, 4. reestimate parameters using the new expectations, 5. iterate until convergence of the likelihood function.

The trick of the EM algorithm thus is to solve the incomplete-data estimation problem for  $\ln g_{\theta}(y)$  indirectly by proceeding iteratively in terms of complete-data estimation for  $\ln p_{\theta}(x)$ . Since the x are not observable,  $\ln p_{\theta}(x)$  is replaced by its conditional expectation given the observed data y and the current fit of the parameter values  $\theta^{(t)}$ . That is, complete-data log-likelihood values are constructed from a conditional expectation given the observed data of the incomplete data problem and the current value of the unknown parameters (E-step). From the thus manufactured complete-data, maximization is simpler and often exists in closed form (M-step). Starting from suitable initial parameter values, the E- and M-steps are iterated until convergence of the incomplete-data log-likelihood L.

More formally, let  $k_{\theta}(x|y) = p_{\theta}(x)/g_{\theta}(y)$  be the conditional probability of x given y and  $\theta$ , then

$$\begin{split} L(\theta') &= \sum_{y \in \mathcal{Y}} \tilde{p}(y) \ln g_{\theta'}(y) \\ &= \sum_{y \in \mathcal{Y}} \tilde{p}(y) k_{\theta} [\ln g_{\theta'}(y)] \\ &= \sum_{y \in \mathcal{Y}} \tilde{p}(y) \sum_{x \in X(y)} k_{\theta}(x|y) \ln \frac{p_{\theta'}(y)}{k_{\theta'}(x|y)} \\ &= \sum_{y \in \mathcal{Y}} \tilde{p}(y) \sum_{x \in X(y)} k_{\theta}(x|y) \ln p_{\theta'}(y) - \sum_{y \in \mathcal{Y}} \tilde{p}(y) \sum_{x \in X(y)} k_{\theta}(x|y) \ln k_{\theta'}(x|y) \\ &= \tilde{p}[k_{\theta}[\ln p_{\theta'}]] - \tilde{p}[k_{\theta}[\ln k_{\theta'}]] \\ &= Q(\theta'; \theta) - H(\theta'; \theta). \end{split}$$

 $Q(\theta'; \theta)$ , the conditional expectation of the complete-data log-likelihood function  $\ln p_{\theta'}(x)$  given y and  $\theta$ , then is used as an auxiliary function to construct an EM algorithm via a mapping  $M : \Theta \to \Theta$ , where each iteration is defined by  $\theta^{(t+1)} = M(\theta^{(t)})$  as follows:

E-step: Compute  $Q(\theta; \theta^{(t)}) = \tilde{p}[k_{\theta^{(t)}}[\ln p_{\theta}]]$ M-step: Choose  $\theta^{(t+1)}$  to be a value of  $\theta \in \Theta$  which maximizes  $Q(\theta; \theta^{(t)})$ .

That is, M is a point-to-set map  $M(\theta^{(t)}) = \underset{\theta \in \Theta}{\arg \max} Q(\theta; \theta^{(t)})$ . This use of Q as an auxiliary function in the EM algorithm can be justified by the fact that an iterative maximization of Q guarantees that the incomplete-data log-likelihood function L is non-decreasing on each iteration of an EM algorithm. This can easily be shown with the inequality

$$L(M(\theta)) - L(\theta) = (Q(M(\theta); \theta) - Q(\theta; \theta)) + (H(\theta; \theta) - H(M(\theta); \theta)) \ge 0, \text{ for all } \theta \in \Theta_{\theta}$$

which follows from the positivity of the difference both in the Q functions (by definition of M) and in the H functions (by Jensen's inequality (see Cover and Thomas (1991))). That is, we have the following proposition, due to Dempster, Laird, and Rubin (1977).

Proposition 4.1 (Dempster et al. (1977), Theorem 1). For each EM algorithm,  $L(M(\theta)) \ge L(\theta)$ , for all  $\theta \in \Theta$ .

Although Q is globally maximized in each M-step, the term H may hinder a straight global maximization of L. As a general result for EM algorithms, Wu (1983) shows that under continuity and differentiability conditions on L and Q, a sequence of EM iterates  $\{L(\theta^{(t)})\}$ bounded from above converges monotonically to a critical point of L.

**Proposition 4.2 (Wu (1983), Theorem 2).** For continuous Q, all limit points of any instance  $\{\theta^{(t)}\}$  of an EM algorithm are critical points of L, and for continuous and differentiable L, a sequence  $\{L(\theta^{(t)})\}$  bounded from above converges monotonically to  $L^* = L(\theta^*)$  for some critical point  $\theta^*$  of L.

To summarize, the popularity of the EM algorithm is due to its easy computation because it relies only on complete-data computations: the E-step involves complete-data conditional expectations, and the M-step requires MLE from these completed data. Even if the algorithm may converge slowly, it conservatively increases the likelihood function at each iteration and in almost all cases converges to a local maximum of L. If a sequence of EM iterates is stuck at some critical point which is not a local or global maximum of L, e.g., a saddle point or even a local minimum, a small random perturbation will help it to diverge from this critical point. If L has several critical points, the convergence properties of an EM sequence will be extremely dependent on the choice of the starting value of the sequence of iterates.

#### 4.3.2 Partial M-Steps: The GEM Algorithm

As discussed in the last section, one main feature of the EM algorithm is to provide a simplified M-step where MLE from complete data rather than from incomplete data is performed. In some cases, even this maximization is complicated and does not exist in closed form. An EM algorithm involving such a complicated M-step would be computationally unattractive. For such cases, Dempster, Laird, and Rubin (1977) defined a so-called generalized EM (GEM) algorithm where the M-step is only partially computed, i.e, each M-step only increases the Q function rather than globally maximizing it.

That is, for a GEM algorithm,  $\theta^{(t+1)}$  is chosen s.t.

$$Q(\theta^{(t+1)}; \theta^{(t)}) \ge Q(\theta^{(t)}; \theta^{(t)})$$

As shown by Dempster, Laird, and Rubin (1977), this condition suffices for increasing the incomplete-data log-likelihood at each interation, i.e., Proposition 4.1 also holds for each GEM

algorithm. However, appropriate convergence of a GEM algorithm does not follow directly without further specification on the process on increasing the Q function. For each instance of a GEM algorithm, one can either show the general convergence conditions for a GEM algorithm as given by Wu (1983) to hold, or directly prove convergence of the specific GEM instance in question. The latter approach is pursued in Sect. 4.6.2 where we explicitly show convergence for a GEM algorithm for log-linear models.

#### 4.3.3 Partial E-steps and Maximum Pseudo-Likelihood Estimation

For many cases, a partial computation of the E-step is also useful. These are especially cases where the sample space  $\mathcal{X}$  is too large to be summed over explicitly in the expectations to be calculated in the E-steps. The idea here is to replace the intractable probability function with respect to which the expectation is taken by a probability function which is more tractable. This change in probability functions results in a corresponding change of the likelihood function to a pseudo-likelihood function which is now defined with respect to the new tractable distribution. Thus from a general optimization-theoretic point of view EM with partial E-steps is an example of maximum pseudo-likelihood estimation.

A theoretical justification for maximum pseudo-likelihood estimation in the context of EM is given in Neal and Hinton (1998) or Csiszár and Tusnády (1984). In terms of Neal and Hinton (1998), the EM algorithm can be seen as maximizing a joint function  $\mathcal{F}$  of the parameters and of the distributions over the unobserved data. Using an arbitrary distribution q over the unobserved variables,  $\mathcal{F}$  can be obtained as a lower bound on the incomplete-data log-likelihood function L as follows.

$$\begin{split} L(\theta) &= \tilde{p}[\ln \sum_{x \in X(y)} p_{\theta}(x)] \\ &= \tilde{p}[\ln \sum_{x \in X(\cdot)} q(x) \frac{p_{\theta}(x)}{q(x)}] \\ &\geq \tilde{p}[\sum_{x \in X(\cdot)} q(x) \ln \frac{p_{\theta}(x)}{q(x)}], \text{ by Jensen's inequality} \\ &= \tilde{p}[q[\ln p_{\theta}]] - \tilde{p}[q[\ln q]] \\ &= \mathcal{F}(q, \theta). \end{split}$$

Provided that values of x are seen as physical states and the energy of a state is  $-\ln p_{\theta}(x)$ , the function  $\mathcal{F}(q, \theta)$  can be seen as analogous to the negative of the "free energy" of statistical physics, i.e., the expected energy under q minus the entropy of q. The EM algorithm can be interpreted in this framework as alternating between maximizing  $\mathcal{F}$  as a function of q and  $\theta$ . The E-step maximizes  $\mathcal{F}$  with respect to q and holds  $\theta$  fixed; the M-step maximizes  $\mathcal{F}$  with respect to  $\theta$  for fixed q.

E-step: Set 
$$q^{(t+1)}$$
 to  $\underset{q}{\operatorname{arg\,max}} \mathcal{F}(q, \theta^{(t)})$ .  
M-step: Set  $\theta^{(t+1)}$  to  $\underset{\theta}{\operatorname{arg\,max}} \mathcal{F}(q^{(t+1)}, \theta)$ 

Neal and Hinton (1998) show that at a true joint maximization, these iterations are equivalent to the classical EM iterations defined in Sect. 4.3.1. That is, the maximum in the Estep is obtained by taking  $q^{(t+1)}(x) = k_{\theta^{(t)}}(x|y)$ , and at this point we have the equality  $\mathcal{F}(q^{(t+1)}, \theta^{(t)}) = L(\theta^{(t)})$ . The maximum in the M-step is obtained by maximizing the term in  $\mathcal{F}$  depending on  $\theta$ , which is in this case  $\tilde{p}[k_{\theta^{(t)}}[\ln p_{\theta}]] = Q(\theta; \theta^{(t)})$ . Since each such E-step guarantees that  $\mathcal{F} = L$ , and since we maximize  $Q(\theta; \theta^{(t)})$  in each M-step, we are guaranteed not to decrease L at each combined EM step.

In a partial E-step,  $q^{(t+1)}$  is set to a tractable approximation of  $k_{\theta^{(t)}}(x|y)$ , which yields the inequality  $\mathcal{F}(q^{(t+1)}, \theta^{(t)}) \leq L(\theta^{(t)})$ . In the corresponding M-step, the term in  $\mathcal{F}$  depending on  $\theta$  is maximized. Together, these combined EM steps guarantee not to decrease the lower bound  $\mathcal{F}$  on the incomplete-data log-likelihood L at each iteration. Thus, for partial E-steps, monotonicity and convergence of the resulting algorithm have to be shown in terms of the pseudo-likelihood function  $\mathcal{F}$  which bounds the true likelihood function L from below.

## 4.4 An EM Example: Baum's Maximization Technique

#### 4.4.1 Basic Concepts

A special instance of the EM algorithm for MLE of hidden Markov models, i.e., stochastic regular grammars, from incomplete data was presented in Baum, Petrie, Soules, and Weiss (1970) and Baum (1972). The form of this algorithm using dynamic programming techniques for efficient computation is well-known as the "forward-backward algorithm" (see Rabiner (1989)). Most popular approaches to parameter estimation for probabilistic grammars are based upon this technique. Baker (1979) generalized this algorithm to the so-called "inside-outside algorithm", which efficiently estimates the parameters of stochastic context-free grammars (see also Booth and Thompson (1973), Lari and Young (1990) and Jelinek, Lafferty, and Mercer (1990)). This algorithm can successfully be applied also to other stochastic grammars which assume independence of their derivation units of each other. Such models are, e.g., stochastic dependency grammars (Carroll and Charniak 1992) or stochastic lexicalised tree-adjoining grammars (Resnik 1992; Schabes 1992). In the following, we will refer to the basic version of this algorithm as Baum's maximization technique.

In the following, we will give a quick review of the basic concepts of Baum's maximization technique. The probabilistic models the algorithm is applied to can be abstracted by stochastic derivation models which define a derivation process as a stochastic process as follows: Make a stochastic choice at each derivation step and assume the stochastic choices to be independent of each other; calculate the probability of a derivation as the joint probability of the independent stochastic choices made, and the probability of an input as the sum of the probabilities of its derivations.

More formally, let  $\pi = (\pi_{ij}) \in \Pi$  be the parameter vector of the probabilistic processing model where  $\pi_{ij} \geq 0$  and  $\sum_j \pi_{ij} = 1$ . The variable *i* ranges over the types of choices that the stochastic process makes, and the variable *j* ranges over the alternatives to choose from when a choice of type *i* is made. Furthermore, let *y* denote an input of the probabilistic processing model, i.e., an observation sequence, and let *x* denote an output of the model, i.e., an analysis, and let Y(x) = y be the unique observation corresponding to analysis *x* and  $X(y) = \{x | Y(x) = y\}$  be the set of analyses of observation *y*. Finally, let  $\nu_{ij}(x)$  be the number of selections of alternative *j* for a choice of type *i* in analysis *x*. The probability of an analysis is the joint probability of the stochastic choices made in producing it. Since these stochastic choices are assumed to be independent of each other, the probability of an analysis is calculated as the product of the probabilities of the stochastic choices made in producing it:

$$p_{\pi}(x) = \prod_{ij} \pi_{ij}^{\nu_{ij}(x)}.$$

The probability of an observation is the sum of the probabilities of its analyses:

$$g_{\pi}(y) = \sum_{x \in X(y)} p_{\pi}(x).$$

For a given random sample of observations, the purpose of Baum's maximization technique is to find maximum likelihood parameter values for the incomplete-data likelihood function Lwhere

$$L(\pi) = \prod_{y \in \mathcal{Y}} g_{\pi}(y)^{\tilde{p}(y)}.$$

The EM mapping M is instantiated here to a particularily simple case. Let  $k_{\pi}(x|y) = p_{\pi}(x)/g_{\pi}(y)$ , then

$$M(\pi_{ij}) = \frac{\tilde{p}[N_{ij}]}{\tilde{p}[\sum_{l} N_{il}]} = \frac{\tilde{p}[k_{\pi}[\nu_{ij}]]}{\tilde{p}[\sum_{l} k_{\pi}[\nu_{il}]]}.$$

Intuitively, the estimated value of parameter  $\pi_{ij}$  is obtained by protating  $N_{ij}$ , the expected number of times choice ij is made during the derivation, by  $\sum_l N_{il}$ , the expected total number

of times a choice of type *i* is made during the derivation, for all observations *y*. Baum, Petrie, Soules, and Weiss (1970) showed that this algorithm is hill-climbing, i.e.,  $L(M(\pi)) \ge L(\pi)$  for all  $\pi \in \Pi$ , and that the incomplete-data likelihood *L* eventually converges to a critical point, i.e., to a local maximum.

#### 4.4.2 Baum's Maximization Technique and Context-Dependence in CLP

The intuitive appeal and the efficient computability of Baum's maximization technique has led to a multiplicity of applications of this technique to various grammar frameworks. Recently, an attempt to apply this technique to a probabilistic version of the constraint-based formalism CUF, which is an instance of the CLP scheme of Höhfeld and Smolka (1988), has been presented by Eisele (1994). As recognized by Eisele (1994), there is a context-dependence problem associated with applying this technique to such constraint-based systems. In CLP terms, the problem is that incompatible variable bindings can lead to failure derivations, which cause a loss of probability mass in the estimated probability distribution over derivations. A similar problem appears in every constraint-based system which constrains derivations by restrictions dependent of the context of the derivation. Approaches embedding Baum's maximization technique into estimation procedures for context-sensitive constraint-based systems have been presented, e.g., by Briscoe and Waegner (1992), Briscoe and Carroll (1993), Brew (1995), Miyata (1996), Osborne and Briscoe (1997) or Bod and Kaplan (1998). From an optimization-theoretic point of view, all such constraint-based approaches contradict the inherent assumptions of Baum's maximization technique which require that the derivation steps are mutually independent and thus the set of licensed derivations is unconstrained.

This problem of context-dependence is discussed in detail in Abney (1997) in connection with the so-called Empirical Relative Frequency (ERF) estimation method, which can be seen a complete-data version of Baum's estimation technique. He shows that applying this method to context-sensitive stochastic attribute-value grammars does not generally yield maximumlikelihood estimates.

In the following, this general argument shall be illustrated with a simple CLP example. Let us apply the stochastic derivation model of Sect. 4.4 to a simple context-sensitive constraint logic program (see Fig. 4.1). The stochastic choices of the abstract model correspond to application probabilities of definite clauses in the generalized SLD-resolution procedure; the alternatives to choose from when an atom is selected in goal reduction are the different clauses defining the selected atom. To indicate a probabilistic parameter  $\pi_{ij}$ , each clause will be annotated by a number ij.

The relational atom  $\mathbf{s}(Z)$  is defined uniquely in clause 11. The atoms  $\mathbf{p}(Z)$  and  $\mathbf{q}(Z)$  each are defined in two different ways, which for the sake of the example are considered to be incompatible. This incompatibility together with the variable sharing in the body of clause 11

 $\mathbf{s}(Z) \leftarrow \mathbf{p}(Z) \& \mathbf{q}(Z).$  $\mathbf{p}(Z) \leftarrow Z = a.$  $\mathbf{p}(Z) \leftarrow Z = b.$  $\mathbf{q}(Z) \leftarrow Z = a.$  $\mathbf{q}(Z) \leftarrow Z = b.$ 

Figure 4.1: Constraint logic program

introduces context-dependence into the program. For a selection of atom p(Z) one can choose between clauses 21 and 22 in a goal reduction step, whereas for a choice of atom q(Z) the alternatives to choose from are clauses 31 and 32.

Suppose we have a training corpus of three queries, consisting of two tokens of query  $y_1 : \mathbf{s}(Z) \& Z = a$  and one token of query  $y_2 : \mathbf{s}(Z) \& Z = b$ . Each query gets a unique proof tree from the program of Fig. 4.1, i.e., a query of type  $y_1$  gets a proof tree of type  $x_1$ , and a query of type  $y_2$  gets one of type  $x_2$  (see Fig. 4.2). Note that in the proof trees of Fig. 4.2 goal reduction and constraint solving are applied in one step.

$$s(Z) \& Z = a | r, c 11, p(Z) \& q(Z) \& Z = a | r, c 21, q(Z) \& Z = a | r, c 31, Z = a$$

$$s(Z) \& Z = b | r, c 11, p(Z) \& q(Z) \& Z = b | r, c 22, q(Z) \& Z = b | r, c 32, Z = b$$

Figure 4.2: Proof trees from constraint logic program

For parameter estimation according to Baum's method, we must calculate conditional probabilities k(x|y) for  $x \in X(y)$ . These probabilities will be 1 in each case, since there is a unique proof tree for each query. Thus for the calculation of  $\tilde{p}[N_{ij}] = \tilde{p}[k_{\pi}[\nu_{ij}]]$ , the expected number of occurences of clauses in proof trees, we simply have to count and can ignore the respective probabilities of the proof trees. As in an application of the complete-data ERF method, for this case Baum's algorithm will give unique parameter estimates  $\hat{\pi}_{ij} = \frac{\tilde{p}[N_{ij}]}{\tilde{p}[\sum_l N_{il}]}$ 

$y\in \mathcal{Y}$	$x \in X(y)$	$ ilde{p}(y)$	k(x y)	$N_{11}$	$N_{21}$	$N_{22}$	$N_{31}$	$N_{32}$
$y_1$	$x_1$	2/3	1	$1 \cdot 1$	$1 \cdot 1$	$1\cdot 0$	$1 \cdot 1$	$1\cdot 0$
$y_2$	$x_2$	1/3	1	$1 \cdot 1$	$1\cdot 0$	$1 \cdot 1$	$1\cdot 0$	$1 \cdot 1$
			$\tilde{p}[N_{ij}] =$	3/3	2/3	1/3	2/3	1/3
			$\tilde{p}[\sum_l N_{il}] =$	3/3	3/3	3/3	3/3	3/3
			$\hat{\pi}_{ij} =$	1	2/3	1/3	2/3	1/3

Table 4.1: Estimation using Baum's maximization technique

in one step (see Table 4.1).

If we consider the calculation of the probability distribution over the proof trees of such a probabilistic CLP model, we see that we cannot simply calculate a product for each proof tree. Instead, we have to introduce a normalization constant in order to ensure the sum over the sample space of proof trees to be 1. For the program of Fig. 4.1, this partition function is taken as the sum of the unnormalized probabilities of the proof trees under the estimated model:  $p_{\hat{\pi}}(x_1) + p_{\hat{\pi}}(x_2) = (1 \cdot 2/3 \cdot 2/3) + (1 \cdot 1/3 \cdot 1/3) = 4/9 + 1/9 = 5/9$ . The normalized probability distribution over proof trees then is:  $p'_{\hat{\pi}}(x_1) = (4/9)/(5/9) = 4/5$ ,  $p'_{\hat{\pi}}(x_2) =$ (1/9)/(5/9) = 1/5. The likelihood L' of our training corpus under the normalized distribution is:  $L' = (4/5)^2 \cdot 1/5 = .128$ . However, note that there is no analytical solution to the problem of finding parameter values  $\pi'$  for the clauses of the program of Fig. 4.1 which define p' as a probabilistic context-free model on the proof trees of Fig. 4.2. Rather, what has happened here is that we implicitly moved to another family of probability distributions by introducing the normalization constant into p'. This new family of probability distributions obviously no longer requires the parameter values to sum up to 1 for identical left-hand sides of rules, but introduces a normalization constant instead in order to guarantee the function to be a probability function. We will acknowledge this family of probability distributions as loglinear distributions in the next section. Clearly, we can easily find parameters of a log-linear model which assigns a higher likelihood to this sample. We could take for example a parameterization  $\pi''$  which assigns  $\pi''_{21} = 2$  and  $\pi''_{ij} = 1$  for all  $ij \neq 21$ . This yields a normalized probability distribution over the proof trees with  $p''_{\pi''}(x_1) = 2/3$ ,  $p''_{\pi''}(x_2) = 1/3$  and likelihood  $L'' = (2/3)^2 \cdot 1/3 = .148$ . The fact that L'' > L' clearly contradicts the assumption that the parameter estimates  $\hat{\pi}$  given by applying Baum's estimation technique to a normalized context-free probability model yield the desired maximum likelihood values.

## 4.5 A Log-Linear Probability Model for CLP

As shown in the last section, we cannot simply apply a stochastic context-free derivation model to CLP but have to go to more expressive probability models. In fact, we implicitly already have made this move in the above example by introducing a partition function into the probabilistic context-free model. We will show in the following that acknowledging this model as a log-linear model not only opens the possibility to find new consistent maximum likelihood estimators but also enables a more flexible parameterization of the probability models.

#### 4.5.1 Motivation

Log-linear models are widely used in probabilistic modelling but come with different names in different applications. The name log-linear is standardly used in contingency table analysis (see, e.g., Knoke and Burke (1980)). The model itself originated under the name of the Gibbs- or Boltzmann-distribution in statistical physics as a flexible probability model of equilibrium states of physical systems. Jaynes (1957) interpreted such equilibrium models in a more abstract framework and coined the name maximum-entropy model. Log-linear models have been applied successfully in the area of image processing, where they are known under the name of random fields (see Geman and Geman (1984)). These special log-linear models are closely related to other probabilistic network models such as Boltzmann machines (see Ackley, Hinton, and Sejnowski (1985)) or Bayesian networks (see Frey (1998)). Loglinear models have been used with effort also in various NLP applications. To name only a few, these applications include probabilistic grammar models (Mark, Miller, Grenander, and Abney 1992; Abney 1997), word spellings (Della Pietra, Della Pietra, and Lafferty 1997), machine translation (Berger, Della Pietra, and Della Pietra 1996), language modelling (Rosenfeld 1996), prepositional phrase attachment (Ratnaparkhi and Roukos 1994), part-of-speech tagging (Ratnaparkhi 1996), history-based parsing (Ratnaparkhi 1997), lexical correlations (Beeferman, Berger, and Lafferty 1997a) text segmentation (Beeferman, Berger, and Lafferty 1997b), and text classification (Nigam, Lafferty, and McCallum 1999).

The popularity of log-linear models is clearly due to the great expressive power they provide with very simple means. That is, log-linear models can be seen as an *exponential* family of probability distributions where the probability of a datum is simply defined as being proportional to the product of weights assigned to selected properties of the datum. Let  $(\pi_i)$ be a vector of weights and  $\nu_i(\omega)$  the number of times property *i* appears in datum  $\omega$ , for all  $i = 1, \ldots, n$ , then

$$p(\omega) \propto \prod_{i=1}^n \pi_i^{
u_i(\omega)}.$$

A log-linear form is obtained from this simply by replacing proportionality by a constant  $C = Z^{\perp 1}$  and parameters  $\pi_i$  by log-parameters  $\lambda_i = \ln \pi_i$ , for all i = 1, ..., n, i.e., taking the logarithm of this probability function yields a linear combination of parameters and properties

and a constant.

$$p(\omega) = C \prod_{i=1}^{n} \pi_{i}^{\nu_{i}(\omega)}$$

$$= Z^{\perp 1} \prod_{i=1}^{n} \pi_{i}^{\nu_{i}(\omega)}$$

$$= Z^{\perp 1} \prod_{i=1}^{n} e^{\lambda_{i}\nu_{i}(\omega)}$$

$$= Z^{\perp 1} e^{\sum_{i=1}^{n} \lambda_{i}\nu_{i}(\omega)}.$$

A more general form of log-linear models is obtained by including a fixed initial or reference distribution  $p_0$  into the model such that  $p(\omega) = Z^{\perp 1} e^{\sum_{i=1}^{n} \lambda_i \nu_i(\omega)} p_0(\omega)$  and  $Z = \sum_{\omega} e^{\sum_i \lambda_i \nu_i(\omega)} p_0(\omega)$ .

Clearly, the main advantage of log-linear models is their great flexibility, which includes the normalized models used in Sect. 4.4.2 and even probabilistic context-free models as special cases (the normalization constant has value 1 in this case). However, considering CLP, with log-linear models we are free to select as properties arbitrary features of proof trees rather than being restricted to clauses only. For example, we could take subtrees of proof trees as properties. This possibility to combine arbitrary clauses to properties allows us to model arbitrary context-dependencies in proof trees. Clearly, linguistically there is no particular reason for assuming rules or clauses as the best properties to use in a probabilistic grammar. As we will see in Sect. 4.7, more abstract properties referring to grammatical functions, attachment preferences, or other general features of constraint-based parses can be employed successfully to probabilistic CLGs. Furthermore, the log-parameters corresponding to these properties are not required to constitute a probability distribution over clauses defining the same predicate, i.e., the parameters do not have to sum to 1 for clauses defining the same predicate. That is, log-linear models allow us to define a probability distribution over proof trees directly rather than indirectly as a joint probability of clause applications as in the context-free models above.

Let us illustrate this with the simple CLP example of Sect. 4.4.2. A training corpus consisting of two tokens of query  $y_1 : \mathbf{s}(Z) \& Z = a$  and one token of query  $y_2 : \mathbf{s}(Z) \& Z = b$  together with the corresponding proof trees generated by the program of Fig. 4.1 is depicted in Fig. 4.3. Note that for ease of readability, we will omit in the following figures the labelings of nodes and edges of proof trees.

To capture the statistics of the training sample of Fig. 4.3, it is sufficient to define a single property which is able to differentiate between the proof tree types. Such a property could be, for example, the terminal node  $\overline{Z=a}$  of proof tree  $x_1$ . Setting the value of the corresponding parameter of this single-parameter model to  $\ln 2$  will yield the desired probability distribution  $p(x_1) = 2/3$ ,  $p(x_2) = 1/3$  with incomplete-data log-likelihood L = .148.

Another way to understand log-linear models is as *maximum-entropy models*. From this viewpoint we do statistical inference and, believing that entropy is the unique consistent measure of the amount of uncertainty represented by a probability distribution, we obey the

$2  imes y_1$ :	$1 imes y_2$ :
$\mathbf{s}(Z) \& Z = a$	$\mathbf{s}(Z) \And Z = b$
$ p(Z) \& q(Z) \& Z = a \\   $	$\begin{array}{c} \mathtt{p}(Z) \And \mathtt{q}(Z) \And Z = b \\ \big  \end{array}$
$\operatorname{q}(Z)\ \&\ Z=a$	$\operatorname{q}(Z) \ \& \ Z = b$
Z = a	Z = b

Figure 4.3: Queries and proof trees for constraint logic program

following principle:

In making inferences on the basis of partial information we must use that probability distribution which has maximum entropy subject to whatever is known. This is the only unbiased assignment we can make; to use any other would amount to arbitrary assumption of information which by hypothesis we do not have. (Jaynes 1957)

More formally, suppose a random variable X can take on values  $x_k, k = 1, \ldots, m$  and we want to estimate the corresponding probabilities  $p_k, k = 1, \ldots, m$ . All we have are expectations of functions  $f_i(X), i = 1, \ldots, n$ . Let these expectations be defined with respect to a given empirical distribution  $\tilde{p}_k, k = 1, \ldots, m$  on complete data  $x_k, k = 1, \ldots, m$  s.t.  $\sum_{k=1}^{m} p_k f_i(x_k) = \sum_{k=1}^{m} \tilde{p}_k f_i(x_k), i = 1, \ldots, n$ . Then the maximum-entropy principle can be stated as follows.

Maximize the entropy  $H(p) = -\sum_{k=1}^{m} p_k$  ln  $p_k$  subject to the constraints  $\sum_{k=1}^{m} p_k f_i(x_k) = \sum_{k=1}^{m} \tilde{p}_k f_i(x_k), i = 1, \dots, n \text{ and } \sum_{k=1}^{m} p_k = 1.$ 

For all  $p_k, k = 1, ..., m$  which solve the above problem, we get the following parametric solution:

$$p_k = \frac{e^{\sum_{i=1}^n \lambda_i f_i(x_k)}}{\sum_{k=1}^m e^{\sum_{i=1}^n \lambda_i f_i(x_k)}}$$

Following Jaynes (1957), this result can be derived directly from a constrained optimization argument where the parameters are viewed as Lagrange multipliers. That is, by applying the standard technique of Lagrange multipliers (see, e.g., Thomas and Finney (1996)) to

the constrained optimization problem stated in the maximum-entropy principle, the above parametric probability model can be derived by solving this constrained optimization problem with respect to the probabilities  $p_k$ . Let  $\Lambda$  denote the Lagrangian defined by

$$\Lambda(p,\lambda) = \sum_{k=1}^{m} (p_k \ln p_k) - (\lambda_0 + 1) (\sum_{k=1}^{m} p_k - 1)$$
$$-\lambda_1 (\sum_{k=1}^{m} p_k f_1(x_k) + \sum_{k=1}^{m} \tilde{p}_k f_1(x_k))$$
$$\vdots$$
$$-\lambda_n (\sum_{k=1}^{m} p_k f_n(x_k) + \sum_{k=1}^{m} \tilde{p}_k f_n(x_k)).$$

Then the first partial derivative of  $\Lambda$  with respect to the  $p_k$  is

$$\frac{\partial}{\partial p_k}\Lambda = (\ln p_k + 1) - (\lambda_0 + 1) - \lambda_1 f_1(x_k) - \dots - \lambda_n f_n(x_k).$$

Now set

$$\frac{\partial}{\partial p_k}\Lambda = 0,$$

then

$$p_k = e^{\lambda_0 + \sum_{i=1}^n \lambda_i f_i(x_k)}.$$

Since the sum of all probabilities  $p_k$  has to be 1, we have

$$1 = \sum_{k=1}^{m} p_k = e^{\lambda_0} \sum_{k=1}^{m} e^{\sum_{i=1}^{n} \lambda_i f_i(x_k)}.$$

If we define a partition function Z as

$$Z = \sum_{k=1}^{m} e^{\sum_{i=1}^{n} \lambda_i f_i(x_k)},$$

then

$$\lambda_0 = \ln Z^{\perp 1}$$

and the maximum-entropy distribution is

$$p_{k} = Z^{\perp 1} e^{\sum_{i=1}^{n} \lambda_{i} f_{i}(x_{k})} = \frac{e^{\sum_{i=1}^{n} \lambda_{i} f_{i}(x_{k})}}{\sum_{k=1}^{m} e^{\sum_{i=1}^{n} \lambda_{i} f_{i}(x_{k})}}.$$

To sum up, the parametric form of maximum-entropy probability models can be derived by solving a constrained optimization problem with respect to the probabilities  $p_k, k = 1, \ldots, m$ .

The remaining problem, namely solving this constrained maximum-entropy problem with respect to the parameters  $\lambda_i$ , i = 1, ..., n, can be shown to be equivalent to solving a maximum likelihood problem for log-linear models. This duality can be stated as follows (see Berger, Della Pietra, and Della Pietra (1996)). The complete-data log-likelihood  $L_c$  of a random sample from a log-linear model  $p_{\lambda}$  on X, with empirical probability  $\tilde{p}(x_k)$  of the values  $x_k, k = 1, ..., m$  is defined as

$$L_c(\lambda) = \ln \prod_{k=1}^m p_\lambda(x_k)^{\tilde{p}(x_k)} = \sum_{k=1}^m \tilde{p}(x_k) \ln p_\lambda(x_k)$$

This function is equivalent to the Lagrangian  $\Lambda$  instantiated to the parametric model  $p_{\lambda}$ :

$$\begin{split} \Lambda(p_{\lambda},\lambda) &= \sum_{k=1}^{m} Z_{\lambda}^{\perp 1} e^{\lambda \cdot f(x_{k})} \ln(Z_{\lambda}^{\perp 1} e^{\lambda \cdot f(x_{k})}) \\ &- \sum_{i=1}^{n} \lambda_{i} \sum_{k=1}^{m} Z_{\lambda}^{\perp 1} e^{\lambda \cdot f(x_{k})} f_{i}(x_{k}) \\ &+ \sum_{i=1}^{n} \lambda_{i} \sum_{k=1}^{m} \tilde{p}(x_{k}) f_{i}(x_{k}) \\ &= -\ln Z_{\lambda} + p_{\lambda} [\lambda \cdot f] - p_{\lambda} [\lambda \cdot f] + \tilde{p} [\lambda \cdot f] \\ &= -\ln Z_{\lambda} + \tilde{p} [\lambda \cdot f]. \end{split}$$

$$L_{c}(\lambda) = \ln \prod_{k=1}^{m} p_{\lambda}(x_{k})^{\tilde{p}(x_{k})}$$
$$= \sum_{k=1}^{m} \tilde{p}(x_{k}) \ln(Z_{\lambda}^{\perp 1} e^{\lambda \cdot f(x_{k})})$$
$$= -\ln Z_{\lambda} + \tilde{p}[\lambda \cdot f].$$

Thus, the values  $\lambda^*$  that solve the constrained maximum-entropy problem with respect to the parameters  $\lambda_i$ , i = 1, ..., n are equivalently a solution to the complete-data maximum likelihood problem for the log-linear model  $p_{\lambda}$ .

The more general model which includes an initial or reference distribution  $p_0$  is derived in a similar way as the unique parametric probability distribution p that minimizes the Kullback Leibler (KL) distance  $D(p||p_0)$  between p and a given reference distribution  $p_0$ , subject to certain constraints. That is, the generalized log-linear model

$$p_{\lambda}(\omega) = \frac{e^{\lambda \cdot \nu(\omega)} p_{0}(\omega)}{\sum_{\omega \in \Omega} e^{\lambda \cdot \nu(\omega)} p_{0}(\omega)}$$

is the parametric solution to the following constraint optimization problem:

Minimize 
$$D(p||p_0) = \sum_{\omega \in \Omega} p(\omega) \ln \frac{p(\omega)}{p_0(\omega)}$$
 subject to the constraints  $p[f_i] = \tilde{p}[f_i], i = 1, ..., n$  and  $\sum_{\omega \in \Omega} p(\omega) = 1$ .

For uniformly distributed  $p_0(\omega)$ , the KL distance  $D(p||p_0)$  is the negative of the entropy H(p), minus a constant not involving  $\lambda$ :

$$D(p||p_0) = \sum_{\omega \in \Omega} p(\omega) \ln p(\omega) - \ln p_0(\omega) = -H(p) - K.$$

In this case, minimizing the KL distance subject to certain constraints is equivalent to maximizing the entropy subject to these constraints. Furthermore, a connection to a maximum likelihood problem can be established for the KL distance minimization problem in a similar way as for the maximum entropy problem.

#### 4.5.2 The Form of Log-Linear Models

Log-linear probability distributions define the probability of a datum simply as proportional to weights assigned to selected properties of the datum. Formally, the parametric family of such distributions is defined as follows.

**Definition 4.1 (Log-linear distribution).** A log-linear probability distribution  $p_{\lambda \cdot \nu}$  on a set  $\Omega$  is defined s.t. for all  $\omega \in \Omega$ :

$$p_{\lambda \cdot \nu}(\omega) = Z_{\lambda \cdot \nu}^{\perp 1} e^{\lambda \cdot \nu(\omega)} p_0(\omega),$$

 $Z_{\lambda \cdot \nu} = \sum_{\omega \in \Omega} e^{\lambda \cdot \nu(\omega)} p_0(\omega) \text{ is a normalizing constant,}$  $\lambda = (\lambda_1, \dots, \lambda_n) \text{ is a vector of log-parameters s.t. } \lambda \in \mathbb{R}^n,$ 

 $\chi = (\chi_1, \ldots, \chi_n)$  is a vector of properties,

 $\nu = (\nu_1, \dots, \nu_n)$  is a vector of property-functions s.t. for each  $\nu_i : \Omega \to \mathbb{N}, \nu_i(\omega)$  is the number of occurences of property  $\chi_i$  in  $\omega$ ,

 $\lambda \cdot \nu(\omega)$  is a weighted property-function s.t.  $\lambda \cdot \nu(\omega) = \sum_{i=1}^{n} \lambda_i \nu_i(\omega)$ ,

 $p_0$  is a fixed initial distribution.

For the following discussion, it will be convenient to introduce some further notation. Properties will be referred to for most purposes by vectors  $\nu$  of property functions rather than by explicit vectors  $\chi$  of properties. Slightly abusing terminology, we will call properties both  $\chi$  and  $\nu$ .

As in Definition 4.1, a log-linear probability distribution depending on property vector  $\nu$ and parameter vector  $\lambda$  will be written in subscript notation as  $p_{\lambda \cdot \nu}$ . In case the property vector is fixed and clear from the context, the model (resp. the normalization constant) will be written  $p_{\lambda}$  (resp.  $Z_{\lambda}$ ) to indicate the dependence on the parameter vector  $\lambda$ .

Furthermore, it will be convenient to have a recursive definition of log-linear models based on weighted property-functions which are extended by additional properties and corresponding parameters.

**Proposition 4.3.** For each weighted property-function  $\phi(\omega) = \lambda \cdot \nu(\omega)$ ,  $\psi(\omega) = \gamma \cdot \mu(\omega)$  let  $(\psi + \phi)(\omega) = \psi(\omega) + \phi(\omega)$  be an extended property-function. Then

$$p_{\psi+\phi}(\omega) = Z_{\psi\circ\phi}^{\perp 1} e^{\psi(\omega)} p_{\phi}(\omega) \text{ where } Z_{\psi\circ\phi} = p_{\phi}[e^{\psi}].$$

Proof.

$$\begin{aligned} p_{\psi+\phi}(\omega) &= Z_{\psi+\phi}^{\perp 1} e^{\psi+\phi(\omega)} p_0(\omega) \\ &= (\sum_{\omega \in \Omega} e^{\psi(\omega)+\phi(\omega)} p_0(\omega))^{\perp 1} e^{\psi(\omega)+\phi(\omega)} p_0(\omega) \\ &= (\sum_{\omega \in \Omega} e^{\psi(\omega)} e^{\phi(\omega)} p_0(\omega) Z_{\phi} Z_{\phi}^{\perp 1})^{\perp 1} e^{\psi(\omega)} e^{\phi(\omega)} p_0(\omega) \\ &= Z_{\phi}^{\perp 1} (\sum_{\omega \in \Omega} e^{\psi(\omega)} p_{\phi}(\omega))^{\perp 1} e^{\psi(\omega)} e^{\phi(\omega)} p_0(\omega) \\ &= (\sum_{\omega \in \Omega} e^{\psi(\omega)} p_{\phi}(\omega))^{\perp 1} e^{\psi(\omega)} p_{\phi}(\omega) \\ &= Z_{\psi\circ\phi}^{\perp 1} e^{\psi(\omega)} p_{\phi}(\omega). \quad \Box \end{aligned}$$

For an extended model with weighted property functions  $\phi(\omega) = \lambda \cdot \nu(\omega)$  and  $\psi(\omega) = \gamma \cdot \nu(\omega)$ , written  $p_{\gamma+\lambda}$ , we have accordingly

$$\begin{split} p_{\gamma+\lambda}(\omega) &= Z_{\gamma+\lambda}^{\perp 1} e^{\gamma \cdot \nu(\omega) + \lambda \cdot \nu(\omega)} p_0(\omega) \\ &= (\sum_{\omega \in \Omega} e^{(\gamma+\lambda) \cdot \nu(\omega)} p_0(\omega))^{\perp 1} e^{(\gamma+\lambda) \cdot \nu(\omega)} p_0(\omega) \\ &= (\sum_{\omega \in \Omega} e^{\gamma \cdot \nu(\omega)} p_{\lambda \cdot \nu}(\omega))^{\perp 1} e^{\gamma \cdot \nu(\omega)} p_{\lambda \cdot \nu}(\omega) \\ &= Z_{\gamma \circ \lambda}^{\perp 1} e^{\gamma \cdot \nu(\omega)} p_{\lambda \cdot \nu}(\omega). \end{split}$$

# 4.6 Statistical Inference for Log-Linear Models from Incomplete Data

In the last two sections we argued that a solution to the context-dependence problem in probabilistic CLP requires probability models which are more expressive than context-free and proposed log-linear models for this purpose. The price we have to pay for this gain in

	log-likelihood	auxiliary function			
complete data	$rac{\partial L_c(\lambda)}{\partial \lambda_i} =  ilde{p}[ u_i] - p_\lambda[ u_i]$	$rac{\partial A_c(\gamma;\lambda)}{\partial \gamma_i} =  ilde{p}[ u_i] - p_{\lambda}[ u_i e^{\gamma_i  u_{\#}}]$			
incomplete data	$\frac{\partial L(\lambda)}{\partial \lambda_i} = \tilde{p}[k_{\lambda}[\nu_i] - p_{\lambda}[\nu_i]]$	$rac{\partial Q(\lambda;\lambda')}{\partial \lambda_i} =  ilde{p}[k_{\lambda'}[ u_i] - p_{\lambda}[ u_i]]$			

Table 4.2: Partial derivatives of objective functions for MLE of log-linear models

expressivity clearly is a gain in complexity of parameter estimation. Furthermore, the gain in flexibility due to property selection is an additional complexity factor which calls for an automatic solution. Fortunately, Della Pietra, Della Pietra, and Lafferty (1997) have presented a statistical inference algorithm for combined property selection and parameter estimation for log-linear models. Abney (1997) has shown the applicability of this algorithm to stochastic attribute-value grammars, which can be seen as a special case of context-sensitive CLGs.

This algorithm, however, applies only to complete data. Unfortunately, the need to rely on large training samples of complete data is a problem if such data are difficult to gather. For example, in natural language parsing applications, complete data means several person-years of hand-annotating large corpora with detailed analyses of specialized grammar frameworks. This is always a labor-intensive and error-prone task, which additionally is restricted to the specific grammar framework, the specific language, and the specific language domain in question. Clearly, for such applications automatic and reusable techniques for statistical inference from incomplete data are desirable.

In the following, we present a version of the statistical inference algorithm of Della Pietra, Della Pietra, and Lafferty (1997) especially designed for incomplete data problems. We present a parameter estimation technique for log-linear models from incomplete data (Sect. 4.6.2) and a property selection procedure from incomplete data (Sect. 4.6.3). These algorithms are combined into a statistical inference algorithm for log-linear models from incomplete data (Sect. 4.6.4). Empirical results on experimenting with these algorithms on a small scale are presented in Sect. 4.7.

This section is based on work presented in shortened form in Riezler (1998a).

#### 4.6.1 Motivation

Why is incomplete-data estimation for log-linear models difficult? The answer is because complete-data estimation for such models is difficult, too. Let us have a look at the first partial derivatives of some objective functions which are considered in MLE of log-linear models from complete and incomplete data (see Table 4.6.1). The system of equations to be solved at the points where the first partial derivatives of the complete data log-likelihood function  $L_c$  are zero, i.e., at the critical points of  $L_c$ , is

$$\sum_{x \in \mathcal{X}} Z_{\lambda}^{\perp 1} e^{\lambda \cdot \nu(x)} \nu_i(x) = \sum_{x \in \mathcal{X}} \tilde{p}(x) \nu_i(x) \text{ for all } i = 1, \dots, n.$$

Clearly, because of the dependence of both  $Z_{\lambda}$  and  $e^{\lambda \cdot \nu(x)}$  on  $\lambda$  this system of equations cannot be solved coordinate-wise in  $\lambda_i$ . This problem is even more severe for the case of incompletedata estimation. The incomplete-data log-likelihood L has its critical points at the solution of the following system of equations in  $\lambda_i$ :

$$\sum_{y \in \mathcal{Y}} \tilde{p}(y) \sum_{x \in \mathcal{X}} Z_{\lambda}^{\perp 1} e^{\lambda \cdot \nu(x)} \nu_i(x) = \sum_{y \in \mathcal{Y}} \tilde{p}(y) \sum_{x \in X(y)} k_{\lambda}(x|y) \nu_i(x) \text{ for all } i = 1, \dots, n.$$

Here additionally a dependence of  $k_{\lambda}(x|y)$  on  $\lambda$  has to be respected. However, an application of the standard EM theory to incomplete-data estimation of log-linear models only partially solves the problem. The equations to be solved to find the critical points of the auxiliary function  $Q(\lambda; \lambda')$  for a log-linear model depending on  $\lambda$  are

$$\sum_{y \in \mathcal{Y}} \tilde{p}(y) \sum_{x \in \mathcal{X}} Z_{\lambda}^{\perp 1} e^{\lambda \cdot \nu(x)} \nu_i(x) = \sum_{y \in \mathcal{Y}} \tilde{p}(y) \sum_{x \in X(y)} k_{\lambda'}(x|y) \nu_i(x) \text{ for all } i = 1, \dots, n.$$

Here  $k_{\lambda'}(x|y)$  depends on  $\lambda'$  instead of  $\lambda$ . However, the dependency of  $Z_{\lambda}$  and  $e^{\lambda \cdot \nu(x)}$  on  $\lambda$  still remains a problem.

Solutions for the system of equations can be found, e.g., by applying general-purpose numerical optimization methods (see Fletcher (1987)) to the problem in question. For the smooth and strictly concave complete-data log-likelihood  $L_c$ , e.g., a conjugate gradient approach could be used. However, optimization methods specifically tailord to the problem of MLE from complete data for log-linear models have been presented by Darroch and Ratcliff (1972) and Della Pietra, Della Pietra, and Lafferty (1997). The "improved iterative scaling" algorithm of Della Pietra, Della Pietra, and Lafferty (1997) itself is an extension of the "generalized iterative scaling" algorithm of Darroch and Ratcliff (1972). In the first algorithm properties are required to sum up to a constant independent of the complete data, i.e.,  $\nu_{\#} = \sum_{i=1}^{n} \nu_i(x) = K$  for all  $x \in \mathcal{X}$ , whereas in the latter algorithm  $\nu_{\#}$  is allowed to vary as a function of x. This property of "generalized iterative scaling" is claimed to improve the convergence rate by increasing the step size taken toward the maximum at each iteration. Both iterative scaling algorithms iteratively maximize an auxiliary function  $A_c(\gamma; \lambda)$  which is defined as a lower bound on the difference  $L_c(\gamma + \lambda) - L_c(\lambda)$  in complete-data log-likelihood when going from a basic model  $p_{\lambda}$  to an extended model  $p_{\gamma+\lambda}$ . The function  $A_c(\gamma; \lambda)$  is maximized as a function of  $\gamma$  for fixed  $\lambda$  which makes it possible to solve the following equation coordinate-wise in  $\gamma_i, i = 1, \ldots, n$ :

$$\sum_{x \in \mathcal{X}} p_{\lambda}(x) \nu_i(x) e^{\gamma_i \nu_{\#}(x)} = \sum_{x \in \mathcal{X}} \tilde{p}(x) \nu_i(x) \text{ for all } i = 1, \dots, n.$$

A closed form solution for  $\gamma_i$  is given for constant  $\nu_{\#}$ ; otherwise simple numerical methods such as Newton's method can be used to solve for the  $\gamma_i$ . It is shown in Della Pietra, Della Pietra, and Lafferty (1997) and Darroch and Ratcliff (1972) that iteratively replacing  $\lambda^{(t+1)}$ by  $\lambda^{(t)} + \gamma^{(t)}$  conservatively increases  $L_c$  and such a sequence of likelihood values eventually converges to the the global maximum of the strictly concave function  $L_c$ .

For the case of incomplete-data estimation things are more complicated. Since the incomplete-data log-likelihood function L is not strictly concave, general-purpose numerical methods such as conjugate gradient cannot be applied. However, such methods can be applied to the auxiliary function Q as defined by a standard EM algorithm for log-linear models. Alternatively, iterative scaling methods can be used to perform maximization of the auxiliary function Q of the EM algorithm. Both approaches result in a doubly iterative algorithm where an iterative algorithm for the M-step is interweaved in the iterative EM algorithm. Clearly, this is computationally burdensome and should be avoided.

The aim of this chapter is exactly to avoid such doubly iterative algorithms. The idea of our approach is to interleave the auxiliary functions Q of the EM algorithm and  $A_c$  of iterative scaling in order to define a singly-iterative incomplete-data estimation algorithm using a new combined auxiliary function. Similar to the case of iterative scaling for complete data, the new auxiliary function will be defined as a lower bound on the improvement in log-likelihood. This allows for an intuitive and elegant proof of convergence of the new algorithm. Our proofs are completely self-contained and do not rely on the convergence of alternating minimization procedures for maximum-entropy models as presented by Csiszár (1975) or Csiszár (1989) or on the regularity conditions for generalized EM algorithms as presented by Wu (1983) or Meng and Rubin (1993). The relation of our algorithm to generalized EM estimation and maximum-entropy estimation is discussed in Sects. 4.6.2.2 and 4.6.2.3.

#### 4.6.2 Parameter Estimation

#### 4.6.2.1 General Theory

Let us start with a problem definition. Applying the incomplete-data framework defined in Sect. 4.3.1 to a log-linear probability model for CLP, we can assume the following to be given:

- observed, incomplete data  $y \in \mathcal{Y}$ , corresponding to a finite sample of queries for a constraint logic program  $\mathcal{P}$ ,
- unobserved, complete data  $x \in \mathcal{X}$ , corresponding to the countably infinite sample of proof trees for queries  $\mathcal{Y}$  from  $\mathcal{P}$ ,
- a many-to-one function  $Y : \mathcal{X} \to \mathcal{Y}$  s.t. Y(x) = y corresponds to the unique query labeling proof tree x, and its inverse  $X : \mathcal{Y} \to 2^{\mathcal{X}}$  s.t.  $X(y) = \{x | Y(x) = y\}$  is the countably infinite set of proof trees for query y from  $\mathcal{P}$ ,

- a complete-data specification  $p_{\lambda}(x)$ , which is a log-linear distribution on  $\mathcal{X}$  with given initial distribution  $p_0$ , fixed property vector  $\chi$  and property-functions vector  $\nu$  and depending on parameter vector  $\lambda$ ,
- an incomplete-data specification  $g_{\lambda}(y)$ , which is related to the complete-data specification by

$$g_{\lambda}(y) = \sum_{x \in X(y)} p_{\lambda}(x)$$

The problem of maximum-likelihood estimation for log-linear models from incomplete data can then be stated as follows.

Given a fixed sample from  $\mathcal{Y}$  and a set  $\Lambda = \{\lambda \mid p_{\lambda}(x) \text{ is a log-linear distribution}$ on  $\mathcal{X}$  with fixed  $p_0$ , fixed  $\nu$  and  $\lambda \in \mathbb{R}^n\}$ , we want to find a maximum likelihood estimate  $\lambda^*$  of  $\lambda$  s.t.  $\lambda^* = \underset{\lambda \in \Lambda}{\operatorname{arg\,max}} L(\lambda) = \ln \prod_{y \in \mathcal{Y}} g_{\lambda}(y)^{\tilde{p}(y)}$ .

For the rest of this section we will refer to a given vector  $\nu$  of property functions. Furthermore, we assume that for each property function  $\nu_i$  some proof tree  $x \in \mathcal{X}$  with  $\nu_i(x) > 0$  exists, and require  $p_{\lambda}$  to be strictly positive on  $\mathcal{X}$ , i.e.,  $p_{\lambda}(x) > 0$  for all  $x \in \mathcal{X}$ . These conditions guarantee that  $p_{\lambda}(x) > 0$  for all  $x \in \mathcal{X}$  and for all  $\lambda \in \Lambda$  which is a desirable property in the following discussion.

Similar to the case of iterative scaling for complete-data estimation, we define an auxiliary function  $A(\gamma, \lambda)$  as a conservative estimate of the difference  $L(\gamma + \lambda) - L(\lambda)$  in log-likelihood. The lower bound for the incomplete-data case can be derived from the complete-data case, in essence, by replacing an expectation of complete, but unobserved data by a conditional expectation given the observed data and the current fit of the parameter values. Clearly, this is the same trick that is used in the EM algorithm, but applied in the context of a different auxiliary function. From the lower-bounding property of the auxiliary function it can immediately be seen that each maximization step of  $A(\gamma, \lambda)$  as a function of  $\gamma$  will increase or hold constant the improvement  $L(\gamma + \lambda) - L(\lambda)$ . This is a first important property of a MLE algorithm. Furthermore, our approach to view the incomplete-data auxiliary function directly as a lower bound on the improvement in incomplete-data log-likelihood enables an intuitive and elegant proof of convergence.

Let the conditional probability of complete data x given incomplete data y and parameter values  $\lambda$  be defined as

$$k_{\lambda}(x|y) = p_{\lambda}(x)/g_{\lambda}(y) = \frac{e^{\lambda \cdot \nu(x)}p_0(x)}{\sum_{x \in X(y)} e^{\lambda \cdot \nu(x)}p_0(x)}.$$

Then a two-place auxiliary function A can be defined as follows.

**Definition 4.2.** Let  $\lambda \in \Lambda$ ,  $\gamma \in \mathbb{R}^n$ ,  $\nu_{\#}(x) = \sum_{i=1}^n \nu_i(x)$ ,  $\bar{\nu}_i(x) = \nu_i(x)/\nu_{\#}(x)$ . Then

$$A(\gamma,\lambda) = \tilde{p}[1 + k_{\lambda}[\gamma \cdot \nu] - p_{\lambda}[\sum_{i=1}^{n} \bar{\nu}_{i}e^{\gamma_{i}\nu_{\#}}]].$$

The particular form of the auxiliary function A and the connection of A and L is discussed in detail in Lemmata 4.5, 4.6, and 4.7 below. Let us first have a look at the extreme value properties of A, which are crucial for the iterative maximization of A.

By considering the first and second derivatives of A, we see that A can be maximized directly and uniquely. This can be explained as follows. Suppose the parameters  $\gamma \in \mathbb{R}^n$  to be a convex set; the Hessian matrix of A is a diagonal matrix filled only with negative elements

$$\frac{\partial^2 A(\gamma;\lambda)}{\partial \gamma_i \partial \gamma_j} = \frac{\partial}{\partial \gamma_j} \left( \frac{\partial A(\gamma;\lambda)}{\partial \gamma_i} \right) = \begin{cases} < 0 & \text{if } i = j \\ 0 & \text{else} \end{cases}$$

and thus negative definite. Unique maximization follows from this since a function whose Hessian is negative definite throughout a convex set is strictly concave, and a strictly concave function attains a maximum at most one point of a convex set, and thus a critical point is necessarily a maximum (see Horn and Johnson (1985)).

**Proposition 4.4.** For each  $\lambda \in \Lambda$ ,  $\gamma \in \mathbb{R}^n$ :  $A(\gamma, \lambda)$  takes its maximum as a function of  $\gamma$  at the unique point  $\hat{\gamma}$  satisfying for each  $\hat{\gamma}_i$ , i = 1, ..., n:

$$\tilde{p}[k_{\lambda}[\nu_i]] = \tilde{p}[p_{\lambda}[\nu_i e^{\hat{\gamma}_i \nu_{\#}}]].$$

Proof.

$$\begin{aligned} \frac{\partial}{\partial \gamma_{i}} A(\gamma, \lambda) &= \frac{\partial}{\partial \gamma_{i}} \tilde{p}[1 + k_{\lambda}[\gamma \cdot \nu] - p_{\lambda}[\sum_{j=1}^{n} \bar{\nu}_{j} e^{\gamma_{j} \nu_{\#}}]] \\ &= \tilde{p}[\frac{\partial}{\partial \gamma_{i}} \sum_{j=1}^{n} (\frac{1}{n} + k_{\lambda}[\gamma_{j} \cdot \nu_{j}] - p_{\lambda}[\bar{\nu}_{j} e^{\gamma_{j} \nu_{\#}}])] \\ &= \tilde{p}[\sum_{j \neq i} (\frac{\partial}{\partial \gamma_{i}} (\frac{1}{n} + k_{\lambda}[\gamma_{j} \cdot \nu_{j}] - p_{\lambda}[\bar{\nu}_{j} e^{\gamma_{j} \nu_{\#}}])) \\ &+ \frac{\partial}{\partial \gamma_{i}} (\frac{1}{n} + k_{\lambda}[\gamma_{i} \cdot \nu_{i}] - p_{\lambda}[\bar{\nu}_{i} e^{\gamma_{i} \nu_{\#}}])] \\ &= \tilde{p}[k_{\lambda}[\nu_{i}] - \sum_{x \in \mathcal{X}} (p_{\lambda}(x)\bar{\nu}_{i}(x)e^{\gamma_{i} \nu_{\#}(x)}\nu_{\#}(x))] \\ &= \tilde{p}[k_{\lambda}[\nu_{i}] - p_{\lambda}[\nu_{i} e^{\gamma_{i} \nu_{\#}}]]. \end{aligned}$$

$$\frac{\partial^2}{\partial \gamma_i^2} A(\gamma, \lambda) = \frac{\partial}{\partial \gamma_i} \tilde{p}[k_\lambda[\nu_i] - p_\lambda[\nu_i e^{\gamma_i \nu_\#}]]$$

$$= -\tilde{p}\left[\frac{\partial}{\partial\gamma_{i}}p_{\lambda}[\nu_{i}e^{\gamma_{i}\nu_{\#}}]\right]$$

$$= -\tilde{p}\left[\sum_{x\in\mathcal{X}}(p_{\lambda}(x)\nu_{i}(x)e^{\gamma_{i}\nu_{\#}(x)}\nu_{\#}(x))\right]$$

$$= -\tilde{p}\left[p_{\lambda}[\nu_{i}\nu_{\#}e^{\gamma_{i}\nu_{\#}}]\right]$$

$$< 0. \square$$

From the auxiliary function A an iterative algorithm for maximizing L is constructed. For want of a name, we will call this algorithm the "Iterative Maximization (IM)" algorithm. At each step of the IM algorithm, a log-linear model based on parameter vector  $\lambda$  is extended to a model based on parameter vector  $\lambda + \hat{\gamma}$ , where  $\hat{\gamma}$  is an estimation of the parameter vector that maximizes the improvement in L when moving away in the parameter space from  $\lambda$ . This increment  $\hat{\gamma}$  is estimated by maximizing the auxiliary function  $A(\gamma, \lambda)$  as a function of  $\gamma$  and, by Proposition 4.4, determined for each  $i = 1, \ldots, n$  uniquely as the solution  $\hat{\gamma}_i$  to the equation  $\tilde{p}[k_{\lambda}[\nu_i]] = \tilde{p}[p_{\lambda}[\nu_i e^{\hat{\gamma}_i \nu_{\#}}]]$ . If  $\nu_{\#} = \sum_{i=1}^n \nu_i(x) = K$  sums to a constant independent of  $x \in \mathcal{X}$ , there exists a closed form solution for the  $\hat{\gamma}_i$ :

$$\hat{\gamma}_i = \frac{1}{K} \ln \frac{\tilde{p}[k_\lambda[\nu_i]]}{p_\lambda[\nu_i]} \text{ for all } i = 1, \dots, n.$$

For  $\nu_{\#}$  varying as a function of x Newton's method can be applied to find an approximate solution (see Sect. 4.8). The IM algorithm in its general form is defined as follows:

**Definition 4.3 (Iterative maximization).** Let  $\mathcal{M} : \Lambda \to \Lambda$  be a mapping defined by

$$\mathcal{M}(\lambda) = \hat{\gamma} + \lambda \text{ with } \hat{\gamma} = \underset{\gamma \in \mathbb{R}^n}{\arg \max} A(\gamma, \lambda)$$

Then each step of the IM algorithm is defined by

$$\lambda^{(k+1)} = \mathcal{M}(\lambda^{(k)}).$$

In order to show the monotonicity and convergence properties of the IM algorithm, we first must prove some provisional results. Lemma 4.5 shows that the auxiliary function  $A(\gamma, \lambda)$  is a lower bound on the incomplete-data log-likelihood difference  $L(\gamma + \lambda) - L(\lambda)$ . In the first inequality we apply Jensen's inequality to the natural logarithm of an expectation. We get a simplified form similar to the log-likelihood difference for complete data, modulo an empirical distribution over complete data being replaced by the conditional distribution  $k_{\lambda}(x|y)$ . This form is simplified further by omitting the logarithm, using the inequality  $\ln x \leq x - 1$ . Furthermore, a random variable  $\nu_{\#}$  on  $\mathcal{X}$  is introduced in order to define a probability distribution  $\bar{\nu}_i$  on  $\mathcal{X}$ . Applying Jensen's inequality to an expectation with respect to  $\bar{\nu}_i$  in the power of e, we arrive at a final simplified form, defining the auxiliary function A.

Lemma 4.5.  $A(\gamma, \lambda) \leq L(\gamma + \lambda) - L(\lambda)$ .

Proof.

$$\begin{split} L(\gamma+\lambda) - L(\lambda) &= \sum_{y \in \mathcal{Y}} \tilde{p}(y) \ln g_{\gamma+\lambda}(y) - \sum_{y \in \mathcal{Y}} \tilde{p} \ln g_{\lambda}(y) \\ &= \tilde{p}[\ln \frac{g_{\gamma+\lambda}(\cdot)}{g_{\lambda}(\cdot)}] \\ &= \tilde{p}[\ln \frac{g_{\gamma+\lambda}(\cdot)}{g_{\lambda}(\cdot)} \sum_{x \in X(\cdot)} (p_{\gamma+\lambda}(x) \frac{p_{\lambda}(x)}{p_{\lambda}(x)})] \\ &= \tilde{p}[\ln \sum_{x \in X(\cdot)} (\frac{p_{\lambda}(x)}{g_{\lambda}(\cdot)} \frac{p_{\gamma+\lambda}(x)}{p_{\lambda}(x)})] \\ &\geq \tilde{p}[\sum_{x \in X(\cdot)} (\frac{p_{\lambda}(x)}{g_{\lambda}(\cdot)} \ln \frac{p_{\gamma+\lambda}(x)}{p_{\lambda}(x)})] \text{ by Jensen's inequality} \\ &= \tilde{p}[\sum_{x \in X(\cdot)} (\frac{p_{\lambda}(x)}{g_{\lambda}(\cdot)} (\ln p_{\gamma+\lambda}(x) - \ln p_{\lambda}(x)))] \\ &= \tilde{p}[\sum_{x \in X(\cdot)} (\frac{p_{\lambda}(x)}{g_{\lambda}(\cdot)} (\ln Z_{\gamma\circ\lambda}^{\perp 1} + \ln e^{\gamma \cdot \nu(x)} + \ln p_{\lambda}(x) - \ln p_{\lambda}(x)))] \\ &= \tilde{p}[k_{\lambda}[\gamma \cdot \nu] - \ln p_{\lambda}[e^{\gamma \cdot \nu}]] \\ &\geq \tilde{p}[k_{\lambda}[\gamma \cdot \nu] + 1 - p_{\lambda}[e^{\gamma \cdot \nu}]] \text{ since } \ln x \leq x - 1 \\ &= \tilde{p}[k_{\lambda}[\gamma \cdot \nu] + 1 - \sum_{x \in \mathcal{X}} (p_{\lambda}(x)e^{\sum_{i=1}^{n} \gamma_{i}\nu_{i}(x)\frac{\nu_{x}(x)}{\nu_{x}(x)}})] \\ &= \tilde{p}[k_{\lambda}[\gamma \cdot \nu] + 1 - \sum_{x \in \mathcal{X}} (p_{\lambda}(x)e^{\sum_{i=1}^{n} \gamma_{i}\nu_{i}(x)\nu_{x}(x)})] \\ &\geq \tilde{p}[k_{\lambda}[\gamma \cdot \nu] + 1 - \sum_{x \in \mathcal{X}} (p_{\lambda}(x)e^{\sum_{i=1}^{n} \gamma_{i}\nu_{i}(x)\nu_{x}(x)})] \\ &\geq \tilde{p}[k_{\lambda}[\gamma \cdot \nu] + 1 - p_{\lambda}[\sum_{i=1}^{n} \bar{\nu}_{i}(x)e^{\gamma_{i}\nu_{x}(x)}] \\ &\geq \tilde{p}[k_{\lambda}[\gamma \cdot \nu] + 1 - p_{\lambda}[\sum_{i=1}^{n} \bar{\nu}_{i}(x)e^{\gamma_{i}\nu_{x}(x)}] \end{bmatrix}$$

Lemma 4.6 shows that there is no estimated improvement in log-likelihood at the origin. Lemma 4.6.  $A(0, \lambda) = 0$ .

Proof.

$$A(0,\lambda) = \tilde{p}[k_{\lambda}[0 \cdot \nu] + 1 - \sum_{x \in \mathcal{X}} p_{\lambda}(x) \sum_{i=1}^{n} \bar{\nu}_i(x) e^0] = 0. \quad \Box$$

Lemma 4.7 shows that the critical points of A and L as functions of  $\gamma$  for fixed  $\lambda$  are the same. Lemma 4.7.  $\frac{d}{dt}\Big|_{t=0} A(t\gamma, \lambda) = \frac{d}{dt}\Big|_{t=0} L(t\gamma + \lambda).$  Proof.

$$\begin{split} \frac{d}{dt}A(t\gamma,\lambda) &= \frac{d}{dt}\bar{p}[k_{\lambda}[t\gamma\cdot\nu] + 1 - \sum_{x\in\mathcal{X}}(p_{\lambda}(x)\sum_{i=1}^{n}\bar{\nu}_{i}(x)e^{t\gamma_{i}\nu_{\#}(x)})] \\ &= \bar{p}[k_{\lambda}[\gamma\cdot\nu] - \sum_{x\in\mathcal{X}}(p_{\lambda}(x)\sum_{i=1}^{n}\frac{\nu_{i}(x)}{\nu_{\#}(x)}e^{t\gamma_{i}\nu_{\#}(x)}\gamma_{i}\nu_{\#}(x))] \\ &= \bar{p}[k_{\lambda}[\gamma\cdot\nu] - \sum_{x\in\mathcal{X}}(p_{\lambda}(x)\sum_{i=1}^{n}\nu_{i}(x)\gamma_{i}e^{t\gamma_{i}\nu_{\#}(x)})]. \\ \frac{d}{dt}\Big|_{t=0}A(t\gamma,\lambda) &= \bar{p}[k_{\lambda}[\gamma\cdot\nu] - \sum_{x\in\mathcal{X}}(p_{\lambda}(x)\sum_{i=1}^{n}\nu_{i}(x)\gamma_{i}e^{0})] \\ &= \bar{p}[k_{\lambda}[\gamma\cdot\nu] - p_{\lambda}[\gamma\cdot\nu]]. \\ \\ \frac{d}{dt}L(t\gamma+\lambda) &= \bar{p}[\frac{d}{dt}\ln\sum_{x\in\mathcal{X})(\cdot)}p_{t\gamma+\lambda}(x)] \\ &= \bar{p}[(\sum_{x\in\mathcal{X})(\cdot)}p_{t\gamma+\lambda}(x))^{\perp 1}\frac{d}{dt}\sum_{x\in\mathcal{X})(\cdot)}e^{t\gamma\cdot\nu(x)}p_{\lambda}(x)Z_{t\gamma\circ\lambda}^{\perp 1}] \\ &= \bar{p}[(\sum_{x\in\mathcal{X})(\cdot)}p_{t\gamma+\lambda}(x))^{\perp 1}\sum_{x\in\mathcal{X})(\cdot)}p_{\lambda}(x)(-e^{t\gamma\cdot\nu(x)}Z_{t\gamma\circ\lambda}^{\perp 2}] \\ &\sum_{x\in\mathcal{X}}e^{t\gamma\cdot\nu(x)}\gamma\cdot\nu(x)p_{\lambda}(x) + Z_{t\gamma\circ\lambda}e^{t\gamma\cdot\nu(x)}\gamma\cdot\nu(x))] \\ &= \bar{p}[-\sum_{x\in\mathcal{X})(\cdot)}p_{t\gamma+\lambda}[\gamma\cdot\nu](\sum_{x\in\mathcal{X})(\cdot)}p_{t\gamma+\lambda}(x))^{\perp 1}] \\ &= \bar{p}[-p_{t\gamma+\lambda}[\gamma\cdot\nu] + k_{t\gamma+\lambda}[\gamma\cdot\nu]]. \\ \\ \\ \\ \frac{d}{dt}\Big|_{t=0}L(t\gamma+\lambda) &= \bar{p}[k_{\lambda}[\gamma\cdot\nu] - p_{\lambda}[\gamma\cdot\nu]]. \\ \end{split}$$

One central result of this section is stated in Theorem 4.8. It shows the monotonicity of the IM algorithm, i.e., the incomplete-data log-likelihood L is increasing on each iteration of the IM algorithm except at fixed points of  $\mathcal{M}$  or equivalently at critical points of L.

**Theorem 4.8 (Monotonicity).** For all  $\lambda \in \Lambda$ :  $L(\mathcal{M}(\lambda)) \geq L(\lambda)$  with equality iff  $\lambda$  is a fixed point of  $\mathcal{M}$  or equivalently is a critical point of L.

Proof.

$$L(\mathcal{M}(\lambda)) - L(\lambda) \geq A(\mathcal{M}(\lambda)) \text{ by Lemma 4.5}$$
  
 
$$\geq 0 \text{ by Lemma 4.6 and definition of } \mathcal{M}.$$

The equality  $L(\mathcal{M}(\lambda)) = L(\lambda)$  holds iff  $\lambda$  is a fixed point of  $\mathcal{M}$ , i.e.,  $\mathcal{M}(\lambda) = \hat{\gamma} + \lambda$  with  $\hat{\gamma} = 0$ . Furthermore,  $\lambda$  is a fixed point of  $\mathcal{M}$  iff  $\hat{\gamma} = \underset{\gamma \in \mathbb{R}^n}{\arg \max} A(\gamma, \lambda) = 0$ ,  $\iff$  for all  $\gamma \in \mathbb{R}^n : \hat{t} = \underset{t \in \mathbb{R}}{\arg \max} A(t\gamma, \lambda) = 0$ ,  $\iff$  for all  $\gamma \in \mathbb{R}^n : \frac{d}{dt}\Big|_{t=0} A(t\gamma, \lambda) = 0$ ,  $\iff$  for all  $\gamma \in \mathbb{R}^n : \frac{d}{dt}\Big|_{t=0} L(t\gamma + \lambda) = 0$ , by Lemma 4.7  $\iff \lambda$  is a critical point of L.

Corollary 4.9 implies that a maximum likelihood estimate is a fixed point of the mapping  $\mathcal{M}$ .

**Corollary 4.9.** Let  $\lambda^* = \underset{\lambda \in \Lambda}{\operatorname{arg max}} L(\lambda)$ . Then  $\lambda^*$  is a fixed point of  $\mathcal{M}$ .

Theorem 4.10 discusses the convergence properties of the IM algorithm. In constrast to the improved iterative scaling algorithm, we cannot show convergence to a global maximum of a strictly concave objective function. Rather we can show convergence of a sequence of IM iterates to a critical point of the non-concave incomplete-data log-likelihood function L. The central property to show is that all limit points of a sequence of IM iterates are critical points of L.

**Theorem 4.10 (Convergence).** Let  $\{\lambda^{(k)}\}$  be a sequence in  $\Lambda$  determined by the IM Algorithm. Then all limit points of  $\{\lambda^{(k)}\}$  are fixed points of  $\mathcal{M}$  or equivalently are critical points of L.

*Proof.* Let  $\{\lambda^{(k_n)}\}$  be a subsequence of  $\{\lambda^{(k)}\}$  converging to  $\overline{\lambda}$ . Then for all  $\gamma \in \mathbb{R}^n$ :

$$\begin{aligned} A(\gamma, \lambda^{(k_n)}) &\leq A(\hat{\gamma}^{(k_n)}, \lambda^{(k_n)}) & \text{by definition of } \mathcal{M} \\ &\leq L(\hat{\gamma}^{(k_n)} + \lambda^{(k_n)}) - L(\lambda^{(k_n)}) & \text{by Lemma 4.5} \\ &= L(\lambda^{(k_n+1)}) - L(\lambda^{(k_n)}) & \text{by definition of IM} \\ &\leq L(\lambda^{(k_n+1)}) - L(\lambda^{(k_n)}) & \text{by monotonicity of } L(\lambda^{(k)}), \end{aligned}$$

and in the limit as  $n \to \infty$ , for continuous A and L:  $A(\gamma, \bar{\lambda}) \leq L(\bar{\lambda}) - L(\bar{\lambda}) = 0$ . Thus  $\gamma = 0$  is a maximum of  $A(\gamma, \bar{\lambda})$ , using Lemma 4.6, and  $\bar{\lambda}$  is a fixed point of  $\mathcal{M}$ . Furthermore,  $\frac{d}{dt}\Big|_{t=0} A(t\gamma, \bar{\lambda}) = \frac{d}{dt}\Big|_{t=0} L(t\gamma + \bar{\lambda}) = 0$ , using Lemma 4.7, and  $\bar{\lambda}$  is a critical point of L.

From this and Theorem 4.8 it follows immediately that each sequence of likelihood values for which an upper bound exists monotonically converges to a critical point of L.

**Corollary 4.11.** Let  $\{L(\lambda^{(k)})\}$  be a sequence of likelihood values bounded from above. Then  $\{L(\lambda^{(k)})\}\$  converges monotonically to a value  $L^* = L(\lambda^*)$  for some critical point  $\lambda^*$  of L.

Thus, the general properties of the IM algorithm are as follows: The IM algorithm conservatively increases the incomplete-data log-likelihood function L. Furthermore, it converges monotonically to a critical point of L, which in almost all cases is a local maximum. And it shows a chaotic behaviour in that for functions L with several extreme values, convergence will be extremely sensitive to the starting value of a sequence of iterates.

#### 4.6.2.2 Relation to Generalized EM Estimation

As discussed in Sect. 4.6.1, a direct application of the standard EM theory to log-linear models is complicated, since complete-data MLE is complicated for log-linear models. That is, a direct application of the EM algorithm to log-linear models always is doubly iterative, because the M-step itself involves some kind of iterative scaling procedure. Examples using iterative M-steps in MLE of log-linear models for partially classified contingency tables are given in Little and Rubin (1987).

Iterative M-steps can be avoided by going to partial M-steps, i.e., to GEM algorithms, as shown in Sect. 4.3.2. In a GEM algorithm, the auxiliary function Q is increased in each M-step rather than maximized. That means, if the improved iterative scaling algorithm is used in the M-step, a single maximization step on the auxiliary function of this algorithm suffices to increase the objective function of this algorithm. Della Pietra, Della Pietra, and Lafferty (1997) use the auxiliary function  $A_c(\gamma, \lambda) = 1 + \tilde{p}[\gamma \cdot \nu] - p_{\lambda}[\sum_{i=1}^{n} \bar{\nu}_i e^{\gamma_i \nu_{\#}}]$  for the objective complete-data log-likelihood function  $L_c(\lambda) = \ln \prod_{x \in \mathcal{X}} p_{\lambda}(x)^{\tilde{p}(x)}$ . An incorporation of this complete-data MLE algorithm into a GEM setting yields the following procedure: First, for a given sample from  $\mathcal{Y}$ , the auxiliary function Q for the incomplete-data log-likelihood  $L = \ln \prod_{y \in \mathcal{Y}} g_{\lambda}(y)^{\tilde{p}(y)}$  is computed as prescribed by the E-step of the EM theory. Next,  $\lambda^{t+1}$ is set to increase Q. That is, we perform only a partial M-step. This task can be fulfilled by tuning the complete-data auxiliary function  $A_c$  of Della Pietra, Della Pietra, and Lafferty (1997) to a new auxiliary function  $\hat{A}$  for the manufactured objective function Q, and by performing a one-step maximization of the complete-data auxiliary function  $A_c$ .

E-step: Compute  $Q(\lambda; \lambda^{(t)}) = \tilde{p}[k_{\lambda^{(t)}}[\ln p_{\lambda}]]$  for a log-linear model  $p_{\lambda}$ .

$$\begin{split} \text{M-step: Choose } \lambda^{(t+1)} \text{ s.t. } Q(\lambda^{(t+1)};\lambda^{(t)}) &\geq Q(\lambda^{(t)};\lambda^{(t)}), \\ \text{i.e., } \lambda^{(t+1)} &= \gamma^{(t)} + \lambda^{(t)} \text{ with } \gamma^{(t)} = \operatorname*{arg\,max}_{\gamma \in \mathbb{R}} \hat{A}(\gamma,\lambda^{(t)}), \\ \text{ and } \hat{A}(\gamma,\lambda^{(t)}) &= \tilde{p}[1+k_{\lambda^{(t)}}[\gamma \cdot \nu] - p_{\lambda^{(t)}}[\sum_{i=1}^{n} \bar{\nu}_{i}e^{\gamma_{i}\nu_{\#}}]]. \end{split}$$

Note that the auxiliary function A which is constructed by applying the complete-data auxiliary function  $A_c$  to the manufactured complete-data log-likelihood Q is identical to our auxiliary function A as specified in Definition 4.2. From the theory of the improved iterative scaling algorithm we can deduce that Q is increased at each M-step of the above procedure. Given this, the theory of the GEM algorithm tells us that the incomplete-data log-likelihood L also is increased at each GEM step of the above procedure. However, convergence of this combined procedure has yet to be studied. An intuitive and elegant way to do this is by considering the auxiliary function A as a lower bound not only on the manufactured complete-data loglikelihood Q but also directly on the incomplete-data log-likelihood L, and prove convergence directly from the relation of A to L. This is the approach we took in the last section.

#### 4.6.2.3 Relation to Maximum-Entropy Estimation

The improved iterative scaling algorithm can be seen also from the perspective of maximumentropy estimation. Della Pietra, Della Pietra, and Lafferty (1997) and Berger, Della Pietra, and Della Pietra (1996) show a duality between maximum likelihood and maximum entropy problems, which can be stated as follows.

The probability distribution  $p^*$  with maximum entropy subject to constraints  $p[f_i] = \tilde{p}[f_i], i = 1, ..., n$  from a distribution  $\tilde{p}(x)$  over complete data  $\mathcal{X}$  is the model in the parametric family of log-linear models  $p_{\lambda}$  that maximizes the likelihood of the training sample  $\mathcal{X}$  distributed according to  $\tilde{p}(x)$ .

Clearly, due to the lack of a distribution  $\tilde{p}(x)$  over complete data  $\mathcal{X}$ , a similar result cannot hold for the incomplete-data case. Rather, in each M-step we get a maximum of a manufactured complete-data likelihood  $Q(\lambda; \lambda') = \tilde{p}[k_{\lambda'}[\ln p_{\lambda}]]$  which corresponds to a maximum-entropy solution subject to constraints from the conditional distribution  $k_{\lambda'}(x|y)$ . If the M-steps are partial themselves , i.e., if we use a GEM setting, then we get the following "increasing-entropy" theorem:

The probability distribution  $p^*$  that increases the entropy H(p) for any probability distribution p subject to the constraints  $p[f_i] = k_{\lambda'}[f_i], i = 1, ..., n$  from a conditional distribution  $k_{\lambda'}(x|y)$  is the model in the parametric family of log-linear probability distributions  $p_{\lambda}$  with  $Q(\lambda; \lambda') \ge Q(\lambda; \lambda)$ .

#### 4.6.3 Property Selection

For the task of parameter estimation discussed in the last section, we assumed a vector of properties to be given. Clearly, exhaustive sets of properties can grow unmanageably large and must be curtailed. An appropriate quality measure on properties can then be used to define an algorithm for automatic property selection.

More generally, property selection can be seen from the viewpoint of model induction. That means, selecting prominent properties out of a set of possible properties can be seen as incrementally inducing a model that captures only the salient statistical qualities of the training data. Such induced models disallow overfitting the training data, which would be the case with models with one unique property per training element. Instead, compact models allow generalizations to new data and temper the overtraining problem.

Different approaches to model induction have been presented. For example, Stolcke and Omohundro (1994) have given a Bayesian approach to inducing the structure of hidden Markov models. This approach starts with a hidden Markov model that directly encodes the data, and proceeds by incrementally generalizing by merging states according to a Bayesian posterior probability measure. This measure trades off the likelihood of the data, which prefers overfitting models, against a prior probability, which prefers simpler models. Maximization of the posterior probability, i.e., the product of the prior and the likelihood, determines which states to merge and when to stop generalizing.

The property selection approach presented by Della Pietra, Della Pietra, and Lafferty (1997) and Berger, Della Pietra, and Della Pietra (1996) proceeds from the opposite direction. Starting from a uniform distribution over the data, which is encoded by a model with no properties at all, properties are incrementally added to the model according to a likelihood measure. A naive form of this measure is the improvement in complete-data log-likelihood when extending a model by a single candidate property c with corresponding log-parameter  $\alpha$ . Unfortunately, when a new parameter is added to the parameter vector of the model, the optimal values can change for all parameters. Thus the calculation of the likelihoodimprovement due to adding a single property requires MLE for all parameters. Clearly, this is infeasible for models with large parameter spaces. Della Pietra, Della Pietra, and Lafferty (1997) and Berger, Della Pietra, and Della Pietra (1996) propose an approximate solution where the complete-data log-likelihood function is maximized directly as a function of a single parameter  $\alpha$ . That is, the improvement due to adding a single candidate is approximated by adjusting only the parameter of this candidate and holding all other parameters fixed. This yields a greedy algorithm which makes it practical to evaluate a large number of candidates at each stage of the combined inference algorithm.

Let us turn now to property selection for log-linear CLP models. For the sake of concreteness, let properties of proof trees be specified as connected, non-overlapping subtrees of proof trees as follows: A property of a proof tree is a connected subgraph of a proof tree, where each node of such a subtree has either zero descendants or the same number of descendants as the corresponding node of the supertree, and the node sets of every two subtrees in the set of properties must not intersect.

Suppose furthermore that properties can be incrementally constructed by selecting from an initial set of goals and from subtrees built by performing a resolution step at a terminal node of a subtree already in the model.

Clearly, an exhaustive set of such properties must be pruned according to some quality measure. What could be an appropriate quality measure for the case of incomplete data? For a MLE framework, the approach of Della Pietra, Della Pietra, and Lafferty (1997) and Berger, Della Pietra, and Della Pietra (1996) offers itself. Unfortunately, we cannot apply the approximate solution of maximizing the likelihood as a function of a single parameter  $\alpha$ , since the incomplete-data log-likelihood L is not concave in the parameters. However, we can express a conservative estimate of the likelihood-gain by instantiating the auxiliary function A of Definition 4.2 to the extension of a model  $p_{\lambda,\nu}$  by a single property c with parameter  $\alpha$ .

$$\begin{aligned} A(\alpha,\lambda) &= \tilde{p}[1+k_{\lambda}[\alpha_{i}c_{i}]-p_{\lambda}[\sum_{i=1}^{n}\bar{c}_{i}e^{\alpha_{i}c_{\#}}]] \\ &= \tilde{p}[1+k_{\lambda}[\alpha c]-p_{\lambda}[e^{\alpha c}]] \\ &\quad \text{since } \alpha_{i}=\alpha, c_{i}(x)=c(x), c_{\#}(x)=c(x), \bar{c}_{i}(x)=1 \end{aligned}$$

From this, we can define an estimated likelihood-gain  $G_c(\alpha, \lambda)$  for a candidate c as follows.

**Definition 4.4.** Let  $\lambda \cdot \nu(x)$  be a weighted property function, c be a candidate property, and  $\alpha \in \mathbb{R}$  the log-parameter corresponding to c. Then the estimated gain  $G_c(\alpha, \lambda)$  of adding candidate property c with parameter value  $\alpha$  to the log-linear model  $p_{\lambda,\nu}$  is defined s.t.

$$G_c(\alpha, \lambda) = \tilde{p}[1 + k_{\lambda \cdot \nu}[\alpha c] - p_{\lambda \cdot \nu}[e^{\alpha c}]].$$

Clearly, this estimated likelihood-gain  $G_c(\alpha, \lambda)$  is a lower bound on the true likelihood-gain  $L(\alpha + \lambda) - L(\lambda)$  for a parameter  $\alpha$  corresponding to a property c.  $G_c(\alpha, \lambda)$  also is strictly concave in the parameters and can be maximized directly and uniquely.

**Proposition 4.12.**  $G_c(\alpha, \lambda)$  takes its maximum as a function of  $\alpha$  at the unique point  $\hat{\alpha}$  satisfying

$$\tilde{p}[k_{\lambda \cdot \nu}[c]] = \tilde{p}[p_{\lambda \cdot \nu}[c \, e^{\hat{\alpha} c}]].$$

Proof.

$$\frac{\partial}{\partial \alpha} G_c(\alpha, \lambda) = \tilde{p}[k_{\lambda \cdot \nu}[c] - p_{\lambda \cdot \nu}[c \ e^{\alpha c}]],$$

$$\frac{\partial^2}{\partial \alpha^2} G_c(\alpha, \lambda) = -\tilde{p}[p_{\lambda \cdot \nu}[c^2 e^{\alpha c}]] < 0. \quad \Box$$

Property selection then will incorporate that property out of the set of candidates that gives the greatest improvement to the model at the property's best adjusted parameter value. Since we are interested only in relative, not absolute gains, a single, non-iterative maximization of the estimated gain will suffice to choose from the candidates. This yields a greedy algorithm for approximate property selection defined as follows. **Definition 4.5 (Property selection).** Let C be a set of candidate properties,  $c \in C$  be a candidate property with log-parameter  $\alpha \in \mathbb{R}$ , and  $G_c(\lambda) = \max_{\alpha} G_c(\alpha, \lambda)$  the maximal estimated gain that property c can give to model  $p_{\lambda \cdot \nu}$ . Then c is selected in a property selection step for model  $p_{\lambda \cdot \nu}$  if  $c = \arg\max_{\alpha} G_{c'}(\lambda)$ .

A reasonable stopping criterion for property selection is to employ cross-validation techniques. That is, the training corpus from  $\mathcal{Y}$  has to be divided into a training portion and a held-out portion. Each candidate property is subjected to maximization of the likelihood for both the training portion and the held-out portion. If the likelihood is increasing for the training portion, but no longer for the held-out portion, the property is discarded. The idea is that at such a point overfitting is indicated for a set of properties that too tighly fits the training portion (and its noise) but no longer provides a good statistical model for both the training and held-out portion of the training corpus. A similar approach of cross-validation can be used to provide a stopping criterion in parameter estimation.

#### 4.6.4 Combined Statistical Inference

The IM procedure for parameter estimation (Definition 4.3) and the procedure for property selection (Definition 4.5) can be combined into a statistical inference algorithm for log-linear models from incomplete data as shown in Table 4.3. The initial model of the Combined Statistical Inference algorithm is assumed to be chosen according to the respective application. For example,  $p_0$  can be chosen as uniform distribution for finite  $\mathcal{X}$ , or as the estimate resulting from an applicaton of Baum's maximization technique to CLP (see Sect. 4.4.2) for infinite  $\mathcal{X}$ . After each property-selection step t, a good starting point for parameter estimation is a  $p_0$  based upon parameter value  $\hat{\alpha} + \lambda^{(t)}$ , where  $\hat{\alpha}$  is the parameter value of the selected property  $\hat{c}$  that maximizes the gain  $G_{\hat{c}}(\alpha, \lambda^{(t)})$ . Note that  $\mathcal{X}$  is defined as the disjoint union of the complete data corresponding to the incomplete data in the random sample, i.e.,  $\mathcal{X} := \sum_{y \in \mathcal{Y} \mid \tilde{p}(y) > 0} X(y)$ .

Let us illustrate this procedure with a simple CLP example. Suppose our sample program is the same as in Fig. 4.1 but with  $\mathcal{L}$ -constraints taken from a language of hierarchical types. The ordering on the types is defined by the operation of set inclusion on the denotations of the types and depicted graphically in Fig. 4.4.



Figure 4.4: Type hierarchy

Furthermore, suppose we have a training corpus of ten queries, consisting of three tokens
**Input** Initial model  $p_0$ , incomplete-data sample from  $\mathcal{Y}$ .

**Output** Log-linear model  $p^*$  on complete-data sample  $\mathcal{X} = \sum_{y \in \mathcal{Y}|\tilde{p}(y)>0} X(y)$ with selected property function vector  $\nu^*$  and log-parameter vector  $\lambda^* = \arg \max_{\lambda \in \Lambda} L(\lambda)$  where  $\Lambda = \{\lambda \mid p_{\lambda} \text{ is a log-linear model on } \mathcal{X} \text{ based on } p_0, \nu^*$ and  $\lambda \in \mathbb{R}^n\}.$ 

### Procedure

- 1.  $p^{(0)} := p_0$  with  $C^{(0)} := \emptyset$ ,
- 2. Property selection: For each candidate property  $c \in C^{(t)}$ , compute the gain  $G_c(\lambda^{(t)}) := \max_{\alpha \in \mathbb{R}} G_c(\alpha, \lambda^{(t)})$ , and select the property  $\hat{c} := \underset{c \in C^{(t)}}{\operatorname{arg}} \max_{c \in C^{(t)}} G_c(\lambda^{(t)})$ .

3. Parameter estimation: Compute a maximum likelihood parameter value  $\hat{\lambda} := \underset{\lambda \in \Lambda}{\operatorname{arg max}} L(\lambda)$  where  $\Lambda = \{\lambda \mid p_{\lambda}(x) \text{ is a log-linear distribution on } \mathcal{X}$  with initial model  $p_0$ , property function vector  $\hat{\nu} := (\nu_1^{(t)}, \nu_2^{(t)}, \ldots, \nu_n^{(t)}, \hat{c}),$  and  $\lambda \in \mathbb{R}^{n+1}\}.$ 

- 4. Until the model converges, set  $p^{(t+1)} := p_{\hat{\lambda} \cdot \hat{\nu}},$  t := t + 1,
  - go to 2.

Table 4.3: Algorithm (Combined Statistical Inference)

of query  $y_1$ :  $\mathbf{s}(Z)$  & Z = a, four tokens of  $y_3$ :  $\mathbf{s}(Z)$  & Z = c, and one token each of query  $y_2$ :  $\mathbf{s}(Z)$  & Z = b,  $y_4$ :  $\mathbf{s}(Z)$  & Z = d, and  $y_5$ :  $\mathbf{s}(Z)$  & Z = e. The corresponding proof trees generated by the program in Fig. 4.1 are given in Fig. 4.5. Note that queries  $y_1, y_2, y_3$  and  $y_4$  are unambiguous, being assigned a single proof tree, while  $y_5$  is ambiguous.

$3  imes y_1$ :	$1 \times y_2$ :	$4 \times y_3$ :
$\mathbf{s}(Z) \And Z = a$	$\begin{array}{c} \mathtt{s}(Z) \And Z = b \\ \big  \end{array}$	$\begin{array}{c} \mathtt{s}(Z) \And Z = c \\ \mid \end{array}$
$ p(Z) \& q(Z) \& Z = a \\   $	$\begin{array}{c} \operatorname{p}(Z) \And \operatorname{q}(Z) \And Z = b \\ \big  \end{array}$	$\begin{array}{c} \mathtt{p}(Z) \And \mathtt{q}(Z) \And Z = c \\ \big  \end{array}$
$\begin{array}{c} \operatorname{q}(Z) \And Z = a \\   \end{array}$	$\begin{array}{c} q(Z) \And Z = b \\ \mid \end{array}$	$\operatorname{q}(Z)\ \&\ Z=a$
Z = a	Z = b	Z = a



Figure 4.5: Queries and proof trees for constraint logic program

A useful first distinction between the proof trees of Fig. 4.5 can be obtained by selecting the two subtrees  $\chi_1 : \mathbb{Z} = a$  and  $\chi_2 : \mathbb{Z} = b$  as properties. These properties allow us to cluster the proof trees in two disjoint sets on the basis of similar statistical qualities of the proof threes in these sets. Since in our training corpus seven out of ten queries come unambiguously with a proof tree including property  $\chi_1$ , we would expect the maximum likelihood parameter value corresponding to property  $\chi_1$  to be higher than the parameter value of property  $\chi_2$ . However, we cannot simply recreate the proportions of the training data from the corresponding proof trees as we did in the unambiguous example of Sect. 4.5. Here we are confronted with an incomplete-data problem, which means that we do not know the frequency of the possible proof trees of query  $y_5$ .

#### 4.7 An Experiment

Let us apply the IM algorithm to this incomplete-data problem. For the selected properties  $\chi_1$  and  $\chi_2$ , we have  $\nu_{\#}(x) = \nu_1(x) + \nu_2(x) = 1$  for all possible proof trees x for the sample of Fig. 4.5. Thus the parameter updates  $\hat{\gamma}_i$  can be calculated from a particularly simple closed form  $\hat{\gamma}_i = \ln \frac{\hat{p}[k_{\lambda}[\nu_i]]}{p_{\lambda}[\nu_i]}$ . A sequence of IM iterates is given in Table 4.4. Probabilities of proof trees involving property  $\chi_i$  are denoted by  $p_i$ . Starting from an initial uniform probability of 1/6 for each proof tree, this sequence of likelihood values converges with an accuracy in the third place after the decimal point after three iterations and yields probabilities  $p_1 \approx .259$  and  $p_2 \approx .074$  for the respective proof trees.

Iteration $t$	$\lambda_1^{(t)}$	$\lambda_2^{(t)}$	$p_1^{(t)}$	$p_2^{(t)}$	$L(\lambda^{(t)})$
0	0	0	1/6	1/6	-17.224448
1	$\ln 1.5$	$\ln .5$	.25	$.08\dot{3}$	-15.772486
2	$\ln 1.55$	$\ln.45$	.2583	.075	-15.753678
3	$\ln 1.555$	$\ln.445$	$.2591\dot{6}$	$.0741\dot{6}$	-15.753481

Table 4.4: Estimation using the IM algorithm

### 4.7 An Experiment

In this section we present an empirical evaluation of the applicability of log-linear probability models and iterative scaling techniques to constraint-based grammars. We present a computationally tractable maximum pseudo-likelihood estimation procedure for log-linear models and apply it to estimating a probabilistic constraint-based grammar from a small corpus of LFG analyses provided by Xerox PARC. The log-linear models employ a small set of about 200 properties to induce a probability distribution on 3000 parses where on average each sentence is ambiguous in 10 parses. The empirical evaluation shows that the correct parse from the set of all parses is found about 59 % of the time.

This section is based on joint work described in Johnson, Geman, Canon, Chi, and Riezler (1999).

### 4.7.1 Incomplete-Data Estimation as Maximum Pseudo-Likelihood Estimation for Complete Data

As we saw in Sect. 4.6, the equations to be solved in statistical inference of log-linear models involve the computation of expectations of property-functions  $\nu_i(x)$  with respect to  $p_{\lambda}(x)$ . Clearly it is possible to find constraint-based grammars where the sample space  $\mathcal{X}$  of parses to be summed over in these expectations is unmanageably large or even infinite.

One possibility to sensibly reduce the summation space is to employ the definition of

the sample space  $\mathcal{X} := \sum_{y \in \mathcal{Y}|\tilde{p}(y)>0} X(y)$  used in incomplete-data estimation as a reduction factor in complete-data estimation. That is, we approximate expectations with respect to the distribution  $p_{\lambda}(\cdot)$  on  $\mathcal{X}$  by considering only such parses  $x \in \mathcal{X}$  whose terminal yield y = Y(x) is seen in the training corpus. Furthermore, the distribution  $g_{\lambda}(y)$  on terminal yields is replaced by the empirical distribution  $\tilde{p}(y)$ :

$$p_{\lambda}[\nu_{i}] = \sum_{x \in \mathcal{X}} p_{\lambda}(x)\nu_{i}(x)$$

$$= \sum_{y \in \mathcal{Y}} \sum_{x \in X(y)} p_{\lambda}(x)\nu_{i}(x)$$

$$= \sum_{y \in \mathcal{Y}} g_{\lambda}(y) \sum_{x \in X(y)} k_{\lambda}(x|y)\nu_{i}(x)$$

$$\approx \sum_{y \in \mathcal{Y}} \tilde{p}(y) \sum_{x \in X(y)} k_{\lambda}(x|y)\nu_{i}(x).$$

Clearly, for most cases the approximate expectation is easier to calculate since the space  $\sum_{y \in \mathcal{V} \mid \tilde{p}(y) > 0} X(y)$  is smaller than the original full space  $\mathcal{X}$ .

The equations to be solved in complete-data estimation for log-linear models are then

$$\sum_{y \in \mathcal{Y}} \tilde{p}(y) \sum_{x \in X(y)} k_{\lambda}(x|y) \nu_i(x) = \sum_{x \in \mathcal{X}} \tilde{p}(x) \nu_i(x) \text{ for all } i = 1, \dots, n.$$

These equations are solutions to the maximization problem of another criterion, namely a complete-data log-pseudo-likelihood function  $PL_c$  which is defined with respect to the conditional probability of parses given the yields observed in the training corpus.

$$PL_c(\lambda) = \ln \prod_{x \in \mathcal{X}, y \in \mathcal{Y}} k_{\lambda}(x|y)^{\tilde{p}(x,y)}$$

In the actual implementation described in Johnson, Geman, Canon, Chi, and Riezler (1999) a slightly different function involving a regularization term promoting small values of  $\lambda$  onto the objective function was maximized. The maximization equations were solved using a conjugategradient approach adapted from Press, Teukolsky, Vetterling, and Flannery (1992). A similar approach to maximum pseudo-likelihood estimation for log-linear models from complete data but in the context of an iterative scaling approach can be found in Berger, Della Pietra, and Della Pietra (1996).

#### 4.7.2 Property Design for Feature-Based CLGs

One central aim of our experiment was to take advantage of the high flexibility of log-linear models and evaluate the usefulness of this issue in hard terms of empirical performance.

The properties employed in our models clearly deviate from the rule or production properties employed in most other probabilistic grammars by encoding as property-functions general

#### 4.7 An Experiment

linguistic principles as proposed by Alshawi and Carter (1994), Srinivas, Doran, and Kulick (1995) or Hobbs and Bear (1995). The definition of properties of LFG parses refers to both the c(onstituent)- and f(eature)-structures of the parses. Examples for the properties employed in our model are

- properties counting the number of adjuncts, arguments and segments in an analysis,
- properties corresponding to grammatical functions used in LFG, including SUBJ, OBJ, OBJ2, COMP, XCOMP, ADJUNCT, etc.
- properties measuring the complexity of the phrase being attached to, thus indicating both high and low attachment,
- properties indicating non-right-branching of nonterminal nodes,
- properties indicating non-parallel coordinate structures,
- properties for atomic attribute-value pairs in feature structures,
- properties for particular syntactic structures such as date-NPs,
- standard rule-properties.

The number of properties defined for each of the two corpora we worked with was about 200 including about 50 rule-properties respectively.

We would also have liked to have included properties corresponding to lexical-semantic head-head relations, but found the small size of our training corpora to be an obstacle in estimating the associated parameters accurately.

### 4.7.3 Empirical Evaluation

The two corpora provided to us by Xeroc PARC contain appointment planning dialogs (Verbmobil corpus, henceforth VM-corpus), and a documentation of Xerox printers (Homecentre corpus, henceforth HC-corpus). The basic properties of the corpora are summarized in Table 4.5. The corpora consist of a packed representation of the c- and f-structures of parses produced for the sentences by a LFG grammar. The LFG parses have been produced automatically by the XLE system (see Maxwell III. and Kaplan (1989)) but corrected manually in addition. Furthermore, it is indicated for each sentence which of its parses is the linguistically correct one. The ambiguity of the sentences in the corpus is 10 parses on average.

In order to cope with the small size of the corpora a 10-way cross-validation framework has been used for estimation and evaluation. That is, the sentences of each corpus were assigned randomly into 10 approximately equal-sized subcorpora. In each run, 9 of the subcorpora

	VM-corpus	HC-corpus
number of sentences	540	980
number of ambiguous sentences	314	481
number of parses of ambiguous sentences	3245	3169

Table 4.5: Properties of the corpora used for the estimation experiment

served as training corpus, and one subcorpus as test corpus. The evaluation scores presented in Tables 4.6 and 4.7 are sums over the the evaluation scores gathered by using each subcorpus in turn as test corpus and training on the 9 remaining subcorpora.

We used two evaluation measures on the test corpus. The first measure  $C_{\text{test}}(\lambda)$  gives the accuracy of disambiguation based on most probable parses. That is,  $C_{\text{test}}(\lambda)$  counts the percentage of sentences in the test corpus whose most probable parse according to a model  $p_{\lambda}$ is the manually determined correct parse. If a sentence has k most probable parses and one of these parses is the correct one, this sentence gets score 1/k. The second evaluation measure is  $-PL_{\text{test}}(\lambda)$ , the negative log-pseudo-likelihood for the correct parses of the test corpus given their yields. This metric measures how much of the probability mass the model puts onto the correct analyses.

In the empirical evaluation, the maximum pseudo-likelihood estimator is compared against a baseline estimator which treats all parses as equally likely. Furthermore, another objective function is considered: The function  $C_{\tilde{\chi}}(\lambda)$  is the number of times the highest weighted parse under  $\lambda$  is the manually determined correct parse in the training corpus  $\tilde{\chi}$ . This function directly encodes the criterion which is used in the linguistic evaluation. However,  $C_{\tilde{\chi}}(\lambda)$  is a highly discontinuous function in  $\lambda$  and hard to maximize. Experiments using a simulated annealing optimization procedure (Press, Teukolsky, Vetterling, and Flannery 1992) for this objective function showed that the computational difficulty of this procedure grows and the quality of the solutions degrades rapidly with the number of properties employed in the model.

The results of the empirical evaluation are shown in Tables 4.6 and 4.7. The maximum pseudo-likelihood estimator performed superior to both the simulated annealing estimator and the uniform baseline estimator on both corpora. The simulated annealing procedure typically scores better than the maximum pseudo-likelihood approach if the number of properties is very small. However, the pseudo-likelihood approach outperforms simulated annealing already for a property-size of 200 as used in our experiment. Furthermore it should be noted that the absolute numbers of 59 % accuracy on the disambiguation task have to be assessed relative to a number of on average 10 parses per sentence.

	$C_{\text{test}}$ for VM-corpus	$-PL_{\text{test}}$ for VM-corpus
uniform baseline estimator	9.7~%	533
simulated annealing estimator	53.7~%	469
maximum pseudo-likelihood estimator	$58.7 \ \%$	396

Table 4.6: Empirical evaluation of estimators on  $C_{\text{test}}$  (accuracy of disambiguation with most probable parse) and  $-PL_{\text{test}}$  (negative log-pseudo-likelihood of correct parses in test corpus) on VM-corpus

	$C_{\text{test}}$ for HC-corpus	$-PL_{\text{test}}$ for HC-corpus
uniform baseline estimator	15.2~%	655
simulated annealing estimator	53.2%	604
maximum pseudo-likelihood estimator	$58.8 \ \%$	583

Table 4.7: Empirical evaluation of estimators on HC-corpus

### 4.8 Approximation Methods

With the algorithms and proofs of the preceding sections in hand, it seems that statistical inference of log-linear models from incomplete data reduces to solving simple equations and computing expectations of simple functions. However, depending on the size of the sample spaces over which these expectations must be taken and depending on the complexity of the parameter- and property-space, these equations can become intractable both analytically and numerically. In order to give a self-contained recipe for statistical inference of log-linear models from incomplete data, we will discuss the possibilities of applying various approximation methods to achieve both analytical and computational tractability in complex applications.

### 4.8.1 Enforcing a Closed-Form Solution

As mentioned above, if the property-functions sum to a constant independent of x, i.e., if

$$\nu_{\#}(x) = \sum_{i=1}^{n} \nu_i(x) = K \text{ for all } x \in \mathcal{X},$$

then the maximum  $\hat{\gamma}$  of the auxiliary function A used in parameter estimation is given in closed form.

For a given vector of property-functions  $\nu$  with  $\nu_{\#}(x) = K$ , the IM algorithm can be stated as shown in Table 4.8. Note that the complete-data sample  $\mathcal{X}$  is computed as  $\mathcal{X} =$ 

 $\sum_{y \in \mathcal{Y} \mid \tilde{p}(y) > 0} X(y).$ 

Input Initial model  $p_0$ , property-functions vector  $\nu$ , incomplete-data sample from  $\mathcal{Y}$ . Output MLE model  $p_{\lambda^*}$  on  $\mathcal{X} = \sum_{y \in \mathcal{Y} \mid \tilde{p}(y) > 0} X(y)$ . Procedure Until convergence do Compute  $p_{\lambda}$ ,  $k_{\lambda}$ , based on  $\lambda = (\lambda_1, \dots, \lambda_n)$ , For *i* from 1 to *n* do  $\gamma_i := \frac{1}{K} \ln \frac{\sum_{y \in \mathcal{Y}} \tilde{p}(y) \sum_{x \in X} (y) k_{\lambda}(x|y) \nu_i(x)}{\sum_{y \in \mathcal{Y}} \sum_{x \in X} p_{\lambda}(x) \nu_i(x)},$   $\lambda_i := \lambda_i + \gamma_i,$ Return  $\lambda^* = (\lambda_1, \dots, \lambda_n)$ .

Table 4.8: Algorithm (Iterative Maximization, Closed-Form)

In this case, the IM algorithm can be seen as an incomplete-data version of the generalized iterative scaling algorithm of Darroch and Ratcliff (1972).

If the constancy-condition is not fulfilled, it can be enforced by introducing a "correction" property-function  $\nu_l$  as follows:

Choose 
$$K = \max_{x \in \mathcal{X}} \nu_{\#}(x)$$
 and  $\nu_l(x) = K - \nu_{\#}(x)$  for all  $x \in \mathcal{X}$   
then  $\sum_{i=1}^{l} \nu_i(x) = K$  for all  $x \in \mathcal{X}$ .

Unfortunately, defining a correction property can be expensive, e.g., in case a property selection procedure is used in statistical inference, a correction property has to be defined after each property selection step.

Correction properties can be avoided by letting  $\nu_{\#}$  vary over  $x \in \mathcal{X}$ . This approach is also claimed to improve the convergence rate of iterative scaling methods by increasing the step size taken toward the maximum at each iteration.

### 4.8.2 Numerical Approximation via Newton's Method

If  $\nu_{\#}(x)$  does not add up to a constant for all  $x \in \mathcal{X}$ , the solutions to the maximization equations in parameter estimation and property selection cannot, in general, be determined in closed form. Fortunately, numerical methods such as Newton's method can be used to efficiently compute approximate solutions to these equations.

#### 4.8 Approximation Methods

Newton's method approximates the solution  $\alpha$  of an equation  $f(\alpha) = 0$  by using a sequence of linearizations of f. At each step, the intersection of the tangent to f at  $\alpha_t$  with the  $\alpha$ -axis is taken, yielding an improved estimate  $\alpha_{t+1}$ . The iteration formulae to approach the solution up to a desired accuracy are defined as follows.

$$\alpha_{t+1} = \alpha_t - \frac{f(\alpha_t)}{f'(\alpha_t)}$$
 where  $f'(\alpha_t)$  is the derivative of  $f$  at  $\alpha_t$ .

This method directly suits our application when we replace  $f(\alpha)$  by the first derivative of the auxiliary function A,  $\frac{\partial}{\partial \gamma_i} A(\gamma, \lambda)$ , in case of parameter estimation, and by the first derivative of the approximate gain  $G_c$ ,  $\frac{\partial}{\partial \alpha} G_c(\alpha, \lambda)$ , in case of property selection. Newton's method usually converges rapidly for such functions.

To efficiently compute the functions in the Newton formulae, we can use a cashing technique similar to the one used in Abney (1997) and apply it to our incomplete-data problem. First, we have to define tables of total probabilities as follows.

- $S_{i,v} = \sum_{x \in \mathcal{X}} p_{\lambda}(x) \delta_{\nu_i(x),v}$  is the expected number of times property function  $\nu_i$  takes value v,
- $T_{i,y} = \sum_{x \in X(y)} k_{\lambda}(x|y) \nu_i(x)$  is the conditionally expected number of times property  $\chi_i$  occurs,
- $U_{i,m} = \sum_{x \in \mathcal{X} \mid \nu_{\#}(x) = m} p_{\lambda}(x) \nu_i(x)$  is the expected number of times property  $\chi_i$  occurs when there is a total number of *m* property instances.

Corresponding to these expectations, we define the following counting variables:

- $s_r(\alpha, i) = \sum_v S_{i,v} e^{\alpha v} v^r$ ,
- $t_r(\alpha, i) = \sum_y T_{i,y} \alpha^r$ ,
- $u_r(\alpha, i) = \sum_m U_{i,m} e^{\alpha m} m^r$ .

The Newton formulae for property selection can then be filled with these expected counts as follows:

$$\begin{aligned} \alpha_{t+1} &= \alpha_t + \frac{\frac{\partial}{\partial \alpha_t} G_c(\alpha_t, \lambda)}{\frac{\partial^2}{\partial \alpha_t^2} G_c(\alpha_t, \lambda)} \\ &= \alpha_t + \frac{\tilde{p}[k_\lambda[c] - Np_\lambda[c \, e^{\alpha_t c}]}{Np_\lambda[c^2 e^{\alpha_t c}]} \\ &= \alpha_t + \frac{t_0(\alpha_t, c) - Ns_1(\alpha_t, c)}{Ns_2(\alpha_t, c)}. \end{aligned}$$

The tables of total probabilities defined above also allow us to express the gain  $G_c(\hat{\alpha}, \lambda)$  of adding property c with best parameter value  $\hat{\alpha}$  to model  $p_{\lambda}$  in terms of expected counts:

$$\begin{aligned} G_c(\hat{\alpha},\lambda) &= N + \tilde{p}[k_\lambda[\hat{\alpha}c] - Np_\lambda[e^{\hat{\alpha}c}] \\ &= N + t_1(\hat{\alpha},c) - Ns_0(\hat{\alpha},c). \end{aligned}$$

For the task of parameter estimation, similar Newton formulae can be obtained from the expected counts:

$$\begin{aligned} \alpha_{t+1} &= \alpha_t + \frac{\frac{\partial}{\partial \alpha_t} A(\gamma, \lambda)}{\frac{\partial^2}{\partial \alpha_t^2} A(\gamma, \lambda)} \\ &= \alpha_t + \frac{\tilde{p}[k_\lambda[\nu_i] - Np_\lambda[\nu_i e^{\alpha_t \nu_\#}]]}{Np_\lambda[\nu_i \nu_\# e^{\alpha_t \nu_\#}]} \\ &= \alpha_t + \frac{t_0(\alpha_t, i) - Nu_0(\alpha_t, i)}{Nu_1(\alpha_t, i)}. \end{aligned}$$

For a random sample from  $\mathcal{Y}$  of size N, an algorithm for approximate parameter estimation can be defined from the above Newton formulae as shown in Table 4.9.

Similarly, an algorithm for approximate property selection can be given as in Table 4.10.

### 4.8.3 Approximating Expectations via Monte Carlo Methods

Independent of whether the solutions of the maximization equations exist in closed form, a further problem arises in connection with large or infinite sample spaces. That is, if the sample space  $\mathcal{X}$  is too large to be summed over in the calculation of the expectations in the maximization equations, methods must be used to approximate these expectations.

One possibility is to use Monte Carlo Methods. Following Abney (1997), we use the Metropolis-Hastings method and show how it can be applied to our incomplete-data problem.

The strategy behind this method is to generate a random sample from a target distribution p by choosing a nominating matrix p' from which sampling is easy, and performing a Bernoulli trial with parameter  $\alpha$  to determine whether to accept or reject the nominated sample point. That means, this method converts a sampler for p' into a sampler for p via an evaluation matrix  $\alpha$ . For our application, we can take as nominating matrix for each query  $y \in \mathcal{Y}$  a stochastic context-free CLP model  $p_{\pi}(x)$  on X(y) as defined in Sect. 4.4.2. From this stochastic derivation model sampling is easy and can be converted by a standard evaluation matrix to sampling from the desired log-linear distribution  $p_{\lambda}(x)$  on X(y).

Input Initial model  $p_0$ , property-functions vector  $\nu$ , incomplete-data sample from  $\mathcal{Y}$ .

**Output** MLE model  $p_{\lambda^*}$  on  $\mathcal{X} = \sum_{y \in \mathcal{Y} | \tilde{p}(y) > 0} X(y)$ .

### Procedure

Until convergence do Compute tables T, U, based on  $\lambda = (\lambda_1, \dots, \lambda_n)$ , For i from 1 to n do  $\alpha := 0$ , Until  $\alpha$  is accurate enough do  $u_0 := 0, u_1 := 0, t_0 := 0$ , For m from 0 to  $m_{\max}$  do  $a := U_{i,m}e^{\alpha m}$ ,  $u_0 := u_0 + a$ ,  $u_1 := u_1 + am$ , For  $y \in \mathcal{Y}$  where  $\tilde{p}(y) > 0$  do  $b := T_{i,y}$ ,  $t_0 := t_0 + b$ ,  $\alpha := \alpha + \frac{t_0 \perp N u_0}{N u_1}$ ,  $\lambda_i := \lambda_i + \alpha$ , Return  $\lambda^* = (\lambda_1, \dots, \lambda_n)$ .

Table 4.9: Algorithm (Iterative Maximization, Newton-Estimate)

```
Input Model p_{\lambda}, set of candidate properties C, incomplete-data sample from \mathcal{Y}.
Output Selected property c^* with maximal parameter value \alpha^*.
Procedure
       Compute tables S, T, based on \lambda,
       G^* := 0, \ c^* := \emptyset, \ \alpha^* := 0,
       For all candidates c \in C do
             \alpha := 0,
             Until \alpha is accurate enough do
                    s_0 := 0, \ s_1 := 0, \ s_2 := 0, \ t_0 := 0, \ t_1 := 0,
                    For v from 0 to v_{\max} do
                            a := S_{c,v} e^{\alpha v},
                            s_0 := s_0 + a,
                            s_1 := s_1 + av,
                            s_2 := s_2 + av^2,
                    For y \in \mathcal{Y} where \tilde{p}(y) > 0 do
                            b := T_{i,y},
                            t_0 := t_0 + b,
                            t_1 := t_1 + b\alpha,
                    \alpha := \alpha + \frac{t_0 \perp N s_1}{N s_2},
             G := N + t_1 - Ns_0,
            If G > G^*, then G^* := G, c^* := c, \alpha^* := \alpha.
       Return c^*, \alpha^*.
```

Table 4.10: Algorithm (Property Selection, Newton-Estimate)

**Input** Initial state  $x_0 \in X(y)$ , Nominating matrix  $p' = p_{\pi}(x)$  on X(y), Log-linear distribution  $p = p_{\lambda}(x)$  on X(y)Evaluation matrix  $\alpha_{x,z} = \begin{cases} 1 & \text{if } p(x)p'(z) \le p(z)p'(x) \\ \frac{p(z)p'(x)}{p(x)p'(z)} & \text{if } p(x)p'(z) > p(z)p'(x) \end{cases}$ Terminal number of steps k**Output** Random sample  $X_0, \ldots, X_k$  from  $p_{\lambda}$  on X(y). Procedure  $X_0 := x_0$ , i := 1, While  $i \leq k$  $x := X_{i\perp 1},$ Randomly generate z from p', If  $z = X_{i\perp 1}$ , then  $X_i := X_{i\perp 1}$ , Else evaluate  $\alpha_{x,z}$ , Randomly generate u from uniform distribution on [0, 1], If  $u \leq \alpha_{x,z}$ , then  $X_i := z$ , Else  $X_i := X_{i \perp 1}$ , i := i + 1, Return  $X_0, \ldots, X_k$ .

Table 4.11: Algorithm (Metropolis-Hastings Sampling)

Following standard textbooks such as Fishman (1996), an application of the Metropolis-Hastings algorithm to our problem is as shown in Table 4.11.

Note that the evaluation matrix  $\alpha_{x,z}$  reduces to a particularly simple form for our application which does not require the computation of normalization constants  $Z_{\lambda}$ . That is, by taking the initial model  $p_0$  of the log-linear CLP model  $p_{\lambda}$  to be of the form of a stochastic CLP model  $p_{\pi}$ , and by assuming independence of the nominated sample points, we get the following form of  $\alpha_{x,z}$ :

$$\alpha_{x,z} = \min\left(1, \frac{p(z)p'(x)}{p(x)p'(z)}\right) \text{ where } \frac{p(z)p'(x)}{p(x)p'(z)} = \frac{p_{\lambda}(z)p_{\pi}(x)}{p_{\lambda}(x)p_{\pi}(z)}$$
$$= \frac{Z_{\lambda}^{\perp 1}e^{\lambda \cdot \nu(z)}p_{\pi}(z)p_{\pi}(x)}{Z_{\lambda}^{\perp 1}e^{\lambda \cdot \nu(x)}p_{\pi}(x)p_{\pi}(z)}$$

$$= \frac{e^{\lambda \cdot \nu(z)}}{e^{\lambda \cdot \nu(x)}}$$
$$= e^{(\lambda \cdot \nu(z) \perp \lambda \cdot \nu(x))}.$$

It can be shown for this sampling method that the distribution of the i.i.d. random variables  $X_i$  converges in distribution to the target distribution  $p_{\lambda}$  as  $i \to \infty$ :

$$\lim_{i \to \infty} P(X_i = x) = p_{\lambda}(x) \text{ for all } x \in X(y).$$

Furthermore, a proper random sample from a probability distribution p enables to estimate expectations of functions f with respect to p directly from the sample points  $X_i$ . That is, the estimated expectation converges to the true expectation with probability 1:

$$\lim_{K \to \infty} \frac{1}{K} \sum_{i=1}^{K} f(X_i) = \sum_{x} f(x) p(x) \text{ with probability 1.}$$

Applying the Metropolis-Hastings algorithm to a log-linear model for CLP yields for each  $y \in \mathcal{Y}$  where  $\tilde{p}(y) > 0$  a random sample  $\tilde{X}(y)$  from  $p_{\lambda}$  on X(y). Such samples can be combined into a sample  $\tilde{\mathcal{X}} = \sum_{y \in \mathcal{Y} | \tilde{p}(y) > 0} \tilde{X}(y)$  from  $p_{\lambda}$  on  $\mathcal{X}$ . From these random samples the desired estimates of expectations of functions with respect to  $p_{\lambda}$  can be computed.

Note that we can use the same random sample for each iteration of Newton's method to estimate the gain for each candidate property simultaneously. After adding the selected property to the model, again a single random sample from the extended model can be used to estimate the MLE values for each parameter in parallel. Suppose we have a random sample from  $\mathcal{Y}$  of size N, a complete data sample  $\tilde{X}(y)$  of size  $M_y$  for y, and combined complete data sample  $\tilde{\mathcal{X}} = \sum_{y \in \mathcal{Y} | \tilde{p}(y) > 0} \tilde{X}(y)$  of size L. Then we can define tables similar to the tables of total probabilities used in Sect. 4.8.2 as follows.

- $S_{i,v} = \sum_{\tilde{x} \in \tilde{\mathcal{X}}} \delta_{\nu_i(\tilde{x}),v}$  is the number of times property function  $\nu_i$  takes value v in combined random sample  $\tilde{\mathcal{X}}$ ,
- $T_{i,y} = \sum_{\tilde{x} \in \tilde{X}(y)} \nu_i(\tilde{x})$  is the number of times property  $\chi_i$  occurs in random sample  $\tilde{X}(y)$ ,
- $U_{i,m} = \sum_{\tilde{x} \in \tilde{\mathcal{X}} | \nu_{\#}(\tilde{x}) = m} \nu_i(\tilde{x})$  is the number of times property  $\chi_i$  occurs in combined random sample  $\tilde{\mathcal{X}}$  when there is a total number of m property instances for each sample point.

For the expectations involved in the closed-form updates in parameter estimation of Sect. 4.8.1, the following counting variables will be convenient:

• 
$$s(i) = \sum_{v} S_{i,v} v$$
,

• 
$$t(i) = \sum_{y} T_{i,y} M_y^{\perp 1}$$
.

The closed-form parameter update  $\hat{\gamma}$  can then be approximated by random sampling as follows.

$$\hat{\gamma}_i \approx \frac{1}{K} \ln \frac{t(i)}{\frac{N}{L}s(i)}.$$

For the expectations involved in the Newton formulae, the counting variables are the same as those of Sect. 4.8.2, except for

• 
$$t_r(\alpha, i) = \sum_y T_{i,y} \alpha^r M_y^{\perp 1}$$
.

The Newton update used in property selection is approximated by random sampling as follows.

$$\alpha_{t+1} \approx \alpha_t + \frac{t_0(\alpha_t, c) - \frac{N}{L}s_1(\alpha_t, c)}{\frac{N}{L}s_2(\alpha_t, c)}.$$

The gain is approximated as

$$G_c(\hat{\alpha}, \lambda) \approx N + t_1(\hat{\alpha}, c) - \frac{N}{L} s_0(\hat{\alpha}, c).$$

Similar random sampling estimates can be obtained for the Newton update used in parameter estimation:

$$\alpha_{t+1} \approx \alpha_t + \frac{t_0(\alpha_t, i) - \frac{N}{L}u_0(\alpha_t, i)}{\frac{N}{L}u_1(\alpha_t, i)}$$

### 4.8.4 Approximating Expectations via Maximum Pseudo-Likelihood Estimation

As stated above, Monte Carlo methods offer the theoretical assurance that the approximation of an expectation converges to the true expectation in the limit. This means that one can get arbitrarily close to the true value of the expectation with increasing sample size. However, convergence can be very slow, i.e., the sample size necessary for an appropriate approximation may be very large. This is especially the case if the distributions of the nominating model  $p_{\pi}$  and the target model  $p_{\lambda}$  are far apart. This may be the case if probabilistic context-free grammars are used as nominating model for a log-linear model on constraint-based grammars. Besides the compensation for sampling errors, many samples may have to be generated to guarantee a reliable estimate of the desired expectations. Together, these problems can make Monte Carlo approximations infeasible in practice. An alternative to Monte Carlo methods is to approximate expectations in a maximum pseudo-likelihood estimation framework. In Sect. 4.3.3 we introduced partial E-steps in the EM algorithm as in instance of maximum pseudo-likelihood estimation. The idea was there to replace an intractable probability function with respect to which an expectation is taken by a probability function which is more tractable. One possibility to achieve such tractable expectations is to use sparse expectations: Instead of replacing the intractable sample space by a Monte-Carlo sample and counting from this, the original sample space is restricted to an appropriate finite subset over which the expectation is calculated.

The general form of such sparse approximations is as follows (cf. Neal and Hinton (1998)). Let S(y) be a finite subset of the set X(y) of complete data corresponding to an incomplete datum  $y \in \mathcal{Y}$ . Then a sparse conditional distribution  $s_{\lambda}(x|y)$  on complete data x given incomplete data y and the current value of the parameters  $\lambda$  can be defined s.t.

$$s_{\lambda^{(t)}}(x|y) = \begin{cases} 0 & \text{if } x \notin S^{(t)}(y), \\ \frac{p_{\lambda^{(t)}}(x)}{\sum_{x \in S^{(t)}(y)} p_{\lambda^{(t)}}(x)} & \text{if } x \in S^{(t)}(y). \end{cases}$$

That is, for a given subset  $S^{(t)}(y)$  of the sample space X(y) defined at time t, the sparse probability distribution  $s_{\lambda^{(t)}}(x|y)$  is defined as the normalized probability distribution that assigns a positive probability only to the elements in  $S^{(t)}(y)$ . The calculation of expectations  $\sum_{x \in S^{(t)}(y)} s_{\lambda^{(t)}}(x|y) f(x)$  of functions f(x) with respect to  $s_{\lambda^{(t)}}(x|y)$  then only takes time proportional to the size of  $S^{(t)}(y)$  at time t.

Various heuristics can be used for a flexible definition of  $S^{(t)}(y)$ . A sensible approach is to define  $S^{(t)}(y)$  as the N most probable  $x \in X(y)$ , and recalculate this set at each step t, and frequently perform a full iteration with  $S^{(t)}(y) = X(y)$  for all  $y \in \mathcal{Y}$  with  $\tilde{p}(y) > 0$ . For N = 1, this approach yields the well-known Viterbi-approximation of the EM algorithm. Here each y is assumed to come with a unique  $x \in X(y)$  at time t. Given algorithms for efficiently searching for the most probable proof tree x for a given query y, a Viterbi-approximation can be defined for parameter estimation of a probabilistic CLP model. A recursive use of this algorithm also enables an N-best-approximation.

A linguistically motivated definition of  $S^{(t)}(y)$  as the trees  $x \in X(y)$  of a context-free grammar which correspond to a bracketing structure annotated to the sample of training sentences has been presented by Pereira and Schabes (1992). Since the bracketing does not change during the estimaton process,  $S^{(t)}(y)$  is constant for all t. Clearly, such bracketing constraints yield on the one hand better linguistic results in terms of a constituent structures of trees consistent with hand-annotated bracketings. On the other hand, the restriction of the sample space to the  $x \in \mathcal{X}$  which correspond to the bracketing structure of the sample from  $\mathcal{Y}$  also reduce the computational load of the estimation process.

A general form of the IM algorithm using sparse approximations  $s_{\lambda}(x|y)$  is given in Table 4.12.

Input Initial model  $p_0$ , initial set  $S_0(y)$ , property-functions vector  $\nu$ , incompletedata sample from  $\mathcal{Y}$ . Output Approximated MLE model  $p_{\lambda^*}$  on  $\mathcal{X} = \sum_{y \in \mathcal{Y} \mid \tilde{p}(y) > 0} X(y)$ . Procedure Until convergence do Compute S(y),  $p_{\lambda}$ ,  $s_{\lambda}$ , based on  $\lambda = (\lambda_1, \dots, \lambda_n)$ , For *i* from 1 to *n* do  $\gamma_i := \frac{1}{K} \ln \frac{\tilde{p}[\sum_{x \in \mathcal{X}} p_{\lambda}(x)\nu_i(x)]}{\sum_{x \in \mathcal{X}} p_{\lambda}(x)\nu_i(x)}$ ,  $\lambda_i := \lambda_i + \gamma_i$ , Return  $\lambda^* = (\lambda_1, \dots, \lambda_n)$ .

Table 4.12: Algorithm (Sparse Iterative Maximization, Closed-Form)

A theoretical justification of such approaches can be given in terms of partial expectations in the context of the EM algorithm. In Sect. 4.3.3, we saw that the incomplete-data loglikelihood  $L(\lambda) = \tilde{p}[\ln g_{\lambda}(y)]$  for a given random sample from  $\mathcal{Y}$  is lower bounded by a pseudolikelihood function  $\mathcal{F}(q, \lambda)$  which is a joint function of the parameters and of the distributions over the unobserved data. The function q can be set to a tractable sparse approximation  $s_{\lambda^{(t)}}$ of  $k_{\lambda^{(t)}}$ . Thus a sparse distribution  $s_{\lambda^{(t)}}$  yields a lower bound  $\mathcal{F}(s_{\lambda^{(t)}}, \lambda^{(t)}) \leq L(\lambda^{(t)})$  in the Estep, which is maximized as a function of  $\lambda$  in the M-step. As shown by Neal and Hinton (1998) or Csiszár and Tusnády (1984), even if some iterations may decrease L, we are guaranteed that the pseudo-likelihood  $\mathcal{F}$  which bounds L from below is increased or held constant with every iteration. The interpretation of the IM algorithm as an instance of a GEM algorithm given in Sect. 4.6.2.2 thus justifies a replacement of  $k_{\lambda^{(t)}}$  by a sparse approximation  $s_{\lambda^{(t)}}$  also for an IM algorithm.

However, it has to be kept in mind that for such partial E-steps monotonicity and convergence of the estimation algorithm has to be proven in terms of the lower bound  $\mathcal{F}$  on L. Clearly, convergence can be shown easily for approaches with constant  $S^{(t)}(y)$  for all t induced, e.g., by fixed bracketing constraints, but is hard to verify for approaches which let  $S^{(t)}(y)$  vary as a function of t such as Viterbi-approximations. More subtle versions of pseudo-likelihood approaches to EM include variational approximation methods, where a parameterized approximating distribution q is used and the parameters are varied to minimize the Kullback Leibler distance between q and  $k_{\lambda}$ . Minimizing this distance clearly results in a minimization of the distance between the pseudo-likelihood function  $\mathcal{F}$  and the true likelihood function L.

$$L(\lambda) - \mathcal{F}(q,\lambda) = ilde{p}[\ln g_{\lambda}(\cdot)] - ilde{p}[\sum_{x \in X(\cdot)} q(x) \ln rac{p_{\lambda}(x)}{q(x)}]$$

$$= \tilde{p}[\ln g_{\lambda}(\cdot) - \sum_{x \in X(\cdot)} q(x)(\ln k_{\lambda}(x|\cdot) + \ln g_{\lambda}(\cdot) - \ln q(x))]$$
  
$$= \tilde{p}[\sum_{x \in X(\cdot)} q(x) \ln \frac{q(x)}{k_{\lambda}(x|\cdot)}]$$
  
$$= \tilde{p}[D(q||k_{\lambda})]$$

The parametric models used, e.g., in the context of large-scale neural networks, are models assuming complete independence of the variables of the network (mean field approximation, see Parisi (1988)) or approximated models probabilistic dependencies of the original model (structured variational approximation, see Saul and Jordan (1996)). Possible applications of variational approximation to estimating probabilistic CLGs could follow these lines. A discussion of such approaches yet is beyond the scope of this thesis.

### 4.9 Parsing and Searching

In the foregoing chapters we discussed the mathematical and algorithmic details of statistical inference of log-linear models from incomplete data, and experimented with these techniques on a small set of real-world data of parses of a constraint-based grammar. On this small scale it was possible to do ambiguity resolution by explicitly listing all parses according to the induced probability distribution and picking the most probable one as the correct one. However, for applications on a larger scale an important question is how the structure of the probability model on parses can be used to guide the search for the most probable parse efficiently without having to list all parses explicitly. Thus the question is whether the search techniques standardly used for probabilistic grammars can be re-applied to the log-linear CLP and CLG models.

We begin our discussion in Sect. 4.9.1 with an application of the tabular parsing method of Earley deduction (Pereira and Warren 1983) to CLGs. The table of pending derivations defined in this method will lay the ground for probabilistic search methods for finding most probable parses. In Sect. 4.9.2 we show that the probabilistic search method of the Viterbi algorithm (Viterbi (1967), Forney (1973)) standardly used in context-free tabular processing models finds the most probable parse of a probabilistic CLG model only under certain restrictions. Since such restrictions may trade off against the search complexity, methods for sensibly relaxing the restrictions are desirable. A heuristic search algorithm resulting from such a relaxation is discussed in Sect. 4.9.3.

### 4.9.1 Earley Deduction for Feature-Based CLGs

Earley deduction has been introduced by Pereira and Warren (1983) as a generalization of Earley's efficient context-free parsing algorithm (Earley (1970), Aho and Ullman (1972)) to a

tabular parsing algorithm for definite clause grammars. In contrast to backtracking methods, in tabular parsing methods a table, or chart, of pending subderivations is built up during derivation. In Earley deduction, subderivations correspond to definite clauses derived from the grammar axioms and a query. Storing such derivation states for future use as items in a chart may avoid the redundancy of backtracking methods which leads in the worst case to an exponential search complexity. Instead, this dynamic-programming technique of storing solutions to subproblems may reduce the search complexity to be polynomial in input length.

The very basic concepts of an application of Earley deduction to CLP can be given as follows. Earley deduction works on two sets of definite clauses, the set of program clauses  $\mathcal{P}$  and the set of derived clauses constituting the chart  $\mathcal{C}$ . An active item of a context-free Earley parser corresponds here to a definite clause with at least one relational atom on its righthandside, i.e., to a non-unit clause. Passive items correspond to clauses whose righthandsides consist only of an  $\mathcal{L}$ -constraint, i.e., to unit-clauses. A selection function determines for each non-unit clause its selected  $\mathcal{R}(\mathcal{L})$ -atom. We adopt here the standard Prolog selection rule where the first atom on the righthandside of a clause is selected in each step. The input to the algorithm consists of a set of program clauses  $\mathcal{P}$  and a query G. The content of the chart  $\mathcal{C}$  initially consists of Gand is continually added to by an exhaustive application of the following two inference rules<sup>2</sup> (the rules are to be read as "If there are clauses  $c_1$  and  $c_2$  and the conditions on these clauses are satisfied, then add clause  $c_3$  to the chart.").

### **Prediction:**

$$c_1 = (H_1 \leftarrow B_1) \in \mathcal{C}$$

$$c_2 = (H_2 \leftarrow B_2) \in \mathcal{P}$$

$$c_3 = (S \leftarrow B'_2 \cup \phi) \in \mathcal{C}$$

where  $c_1$  is non-unit,  $c_2$  is unit or non-unit, S is the selected atom in  $B_1$ ,  $\phi$  is the  $\mathcal{L}$ -constraint in  $B_1$ , and there exists a variant  $c'_2 = (S \leftarrow B'_2)$  of  $c_2$  s.t.  $V(c_1) \cap V(B'_2) \subseteq V(S)$  and the  $\mathcal{L}$ -constraint  $\phi'$  of  $c_3$  is satisfiable.

<sup>&</sup>lt;sup>2</sup>Prediction is called "instantiation" in Pereira and Warren (1983) and completion corresponds to their "resolution". In context-free Earley parsing standardly a distinction between "predictor", "scanner" and "completer" operations is made. The first operation corresponds to prediction and the latter two operations are subsumed by the completion operation of the Earley decuction framework defined below.

#### **Completion:**

$$c_1 = (H_1 \leftarrow B_1) \in \mathcal{C}$$

$$c_2 = (H_2 \leftarrow B_2) \in \mathcal{C}$$

$$c_3 = (H_1 \leftarrow (B_1 \setminus S) \cup B'_2) \in \mathcal{C}$$

where  $c_1$  is non-unit,  $c_2$  is unit, S is the selected atom in  $B_1$ , and there exists a variant  $c'_2 = (S \leftarrow B'_2)$  of  $c_2$  s.t.  $V(c_1) \cap V(B'_2) \subseteq V(S)$  and the  $\mathcal{L}$ -constraint  $\phi'$  of  $c_3$  is satisfiable.

These rules can be rationalized as follows: The prediction rule proposes for the selected atom of a clause  $c_1$  a possible variant of a program clause  $c_2$  using which an  $\xrightarrow{r,c}$ -step, i.e., a combined goal-reduction and constraint-solving step, can be performed. For a unit clause  $c_2$ , the completion rule then performs a combined  $\xrightarrow{r,c}$ -step on the lefthandside atom of  $c_2$  and substitutes this selected atom in clause  $c_1$  by the resulting righthandside  $\mathcal{L}$ -constraint. Both rules collect the  $\mathcal{L}$ -constraints of the antecedent clauses and take care of successful constraint solving and prevent accidental variable sharing in the consequent clause.

Clearly, this combination of top-down prediction and bottom-up completion defines a search rule which can reduce the parsing complexity in comparison to backtracking methods. However, to make these inference rules a workable algorithm, several issues concerning the effective applicability of Earley deduction to different purposes have to be addressed. Since these topics are not of direct relevance for our problem, we refer the reader to the extensive literature on this subject (see, e.g., Pereira and Shieber (1987), Dörre (1993), Dörre and Johnson (1995)).

Let us illustrate the basic concepts of Earley deduction with a simple feature-based CLG. In the following example we will make use of a standard technique for string position indexing, e.g., the indexed clause

$$\operatorname{sign}(X, 0, 1) \leftarrow X = \phi.$$

abbreviates an actual CLP clause

$$\operatorname{sign}(X, Y, Z) \leftarrow X = \phi \& Y = 0 \& Z = 1.$$

where the constants 0 and 1 denote the start and end position of the span of the predicate in the input string. The string position can be read off for unit clauses from the lefthandside atom, but for non-unit clauses from the first string position argument of the head atom and the first string position argument of the leftmost atom in the body. Note that string position indexing is not mentioned in the definition of the inference rules for Earley deduction. In fact, this indexing is not necessary for Earley deduction to work. Rather, it is an effective way to reduce the number of unsuccesful rule applications in an implementation, and will also make our example more transparent.

Let us return for illustration to the simple example of Fig. 3.5. An indexed variant of this program is given in Fig. 4.6.

1 phrase $(X, S_0, S) \leftarrow X = (phrase \land CAT : s \land DTR1:CAT : n \land DTR2:CAT : v \land DTR1:AGR : Y \land DTR2:AGR : Y \land DTR1 : Z_1 \land DTR2 : Z_2) \& \operatorname{sign}(Z_1, S_0, S_1) \& \operatorname{sign}(Z_2, S_1, S).$ 

2 phrase $(X, S_0, S) \leftarrow X = (phrase \wedge CAT : np \wedge DTR1: CAT : n \wedge DTR2: CAT : n \wedge DTR1 : Z_1 \wedge DTR2 : Z_2) \& sign(Z_1, S_0, S_1) \& sign(Z_2, S_1, S).$ 

 $3 \operatorname{word}(X, 0, 1) \leftarrow X = (word \wedge \operatorname{CAT} : n \wedge \operatorname{PHON} : Clinton \wedge \operatorname{AGR} : sg).$ 

4 word $(X, 1, 2) \leftarrow X = (word \land CAT : v \land PHON : talks \land AGR : sg).$ 

5 word $(X, 1, 2) \leftarrow X = (word \land CAT : n \land PHON : talks \land AGR : pl).$ 

 $6 \operatorname{sign}(X, S_0, S) \leftarrow \operatorname{phrase}(X, S_0, S).$ 

 $7 \operatorname{sign}(X, S_0, S) \leftarrow \operatorname{word}(X, S_0, S).$ 

Figure 4.6: Indexed feature-based constraint logic grammar

An application of Earley deduction to parsing the query

 $sign(X, 0, 2) \& X = (sign \land DTR1: PHON : Clinton \land DTR2: PHON : talks).$ 

denoting the input sentence

#### $_0Clinton_1 talks_2.$

is given in Figs. 4.7 and 4.8.

A convenient way to illustrate graphically the relation of derived items in a chart to partial parses of an input sentence is by a chart graph. An chart graph for the sequence of derived clauses of Figs. 4.7 and 4.8 is given in Fig. 4.9. This graph associates the numbers of derived clauses with directed edges which connect input string position nodes. Edges which point from a node to the node itself are attached with the numbers of clauses derived by prediction with non-unit clauses. In the language of context-free Earley parsing, such edges represent predictions on non-terminal symbols. Scanning of terminal symbols is represented by edges connecting a node with the next node on its right. Such edges are attached with the numbers of clauses derived by prediction with unit clauses.

- 9  $\leftarrow \operatorname{sign}(X, 0, 2) \& X = (sign \land \operatorname{DTR1}: \operatorname{PHON} : Clinton \land \operatorname{DTR2}: \operatorname{PHON} : talks).$  (I)
- 10  $\operatorname{sign}(X, 0, 2) \leftarrow \operatorname{phrase}(X, 0, 2) \& X = (sign \land DTR1: PHON : Clinton \land (P 9, 6) DTR2: PHON : talks).$
- 11 phrase $(X, 0, 2) \leftarrow \operatorname{sign}(Z_1, 0, S_1)$  & sign $(Z_2, S_1, 2)$  &  $X = (phrase \land \operatorname{CAT} : s \land \operatorname{DTR1} : (P 10, 1)$ word  $\land$  DTR1: CAT :  $n \land$  DTR1: PHON : Clinton  $\land$  DTR1: AGR :  $Y \land$  DTR2 : word  $\land$  DTR2: CAT :  $v \land$  DTR2: PHON : talks  $\land$  DTR2: AGR :  $Y \land$  DTR1 :  $Z_1 \land$  DTR2 :  $Z_2$ ).
- 12  $\operatorname{sign}(Z_1, 0, S_1) \leftarrow \operatorname{word}(Z_1, 0, S_1) \& X = (phrase \land CAT : s \land DTR1 : word \land (P 11,7)$ DTR1: CAT :  $n \land DTR1$ : PHON : Clinton  $\land DTR1$ : AGR :  $Y \land DTR2$  : word  $\land DTR2$ : CAT :  $v \land DTR2$ : PHON : talks  $\land DTR2$ : AGR :  $Y \land DTR1$  :  $Z_1 \land DTR2 : Z_2$ ).
- 13 word $(Z_1, 0, 1) \leftarrow X = (phrase \land CAT : s \land DTR1 : word \land DTR1 : CAT : (P 12, 3)$   $n \land DTR1 : PHON : Clinton \land DTR1 : AGR : Y \land DTR1 : AGR : sg \land DTR2 :$   $word \land DTR2 : CAT : v \land DTR2 : PHON : talks \land DTR2 : AGR : Y \land DTR2 : AGR :$  $sg \land DTR1 : Z_1 \land DTR2 : Z_2).$
- 14 phrase $(X, 0, 2) \leftarrow \operatorname{sign}(Z_1, 0, S_1)$  & sign $(Z_2, S_1, 2)$  &  $X = (phrase \land \operatorname{CAT} : np \land (P 10, 2)$ DTR1: word  $\land$  DTR1: CAT :  $n \land$  DTR1: PHON : Clinton  $\land$  DTR2 : word  $\land$ DTR2: CAT :  $n \land$  DTR2: PHON : talks  $\land$  DTR1 :  $Z_1 \land$  DTR2 :  $Z_2$ ).
- 15  $\operatorname{sign}(Z_1, 0, S_1) \leftarrow \operatorname{word}(Z_1, 0, S_1) \& X = (phrase \land \operatorname{CAT} : np \land \operatorname{DTR1} : word \land (P 14, 7)$ DTR1: CAT :  $n \land$ DTR1: PHON : Clinton  $\land$ DTR2 : word  $\land$ DTR2: CAT :  $n \land$ DTR2: PHON : talks  $\land$ DTR1 :  $Z_1 \land$ DTR2 :  $Z_2$ ).
- 16 word $(Z_1, 0, 1) \leftarrow X = (phrase \land CAT : np \land DTR1 : word \land DTR1 : CAT : (P 15,3)$  $<math>n \land DTR1 : PHON : Clinton \land DTR1 : AGR : sg \land DTR2 : word \land DTR2 : CAT :$  $n \land DTR2 : PHON : talks \land DTR1 : Z_1 \land DTR2 : Z_2).$
- 17  $\operatorname{sign}(Z_1, 0, 1) \leftarrow X = (phrase \land CAT : s \land DTR1 : word \land DTR1 : CAT : (C 12, 13)$   $n \land DTR1 : PHON : Clinton \land DTR1 : AGR : Y \land DTR1 : AGR : sg \land DTR2 :$   $word \land DTR2 : CAT : v \land DTR2 : PHON : talks \land DTR2 : AGR : Y \land DTR2 : AGR :$  $sg \land DTR1 : Z_1 \land DTR2 : Z_2).$
- 18 phrase $(X, 0, 2) \leftarrow sign(Z_2, 1, 2) \& X = (phrase \land CAT : s \land DTR1 : word \land (C 11, 17)$ DTR1: CAT :  $n \land DTR1$ : PHON : Clinton  $\land DTR1$ : AGR :  $Y \land DTR1$ : AGR :  $sg \land DTR2 : word \land DTR2$ : CAT :  $v \land DTR2$ : PHON : talks  $\land DTR2$ : AGR :  $Y \land DTR2$ : AGR :  $sg \land DTR1 : Z_1 \land DTR2 : Z_2$ ).
- **19**  $\operatorname{sign}(Z_1, 0, 1) \leftarrow X = (phrase \land CAT : np \land DTR1 : word \land DTR1 : CAT : (C 15, 16)$  $n \land DTR1 : PHON : Clinton \land DTR1 : AGR : sg \land DTR2 : word \land DTR2 : CAT :$  $n \land DTR2 : PHON : talks \land DTR1 : Z_1 \land DTR2 : Z_2).$
- 20 phrase $(X, 0, 2) \leftarrow \operatorname{sign}(Z_2, 1, 2)$  &  $X = (phrase \land CAT : np \land DTR1 : word \land (C 14, 19)$ DTR1: CAT :  $n \land DTR1$ : PHON : Clinton  $\land DTR1$ : AGR :  $sg \land DTR2$  : word  $\land DTR2$ : CAT :  $n \land DTR2$ : PHON : talks  $\land DTR1 : Z_1 \land DTR2 : Z_2$ ).

Figure 4.7: Earley deduction chart

- 21  $\operatorname{sign}(Z_2, 1, 2) \leftarrow \operatorname{word}(Z_2, 1, 2) \& X = (phrase \land \operatorname{CAT} : s \land \operatorname{DTR1} : word \land (P 18, 7)$ DTR1: CAT :  $n \land \operatorname{DTR1}$ : PHON : Clinton  $\land \operatorname{DTR1}$ : AGR :  $Y \land \operatorname{DTR1}$ : AGR :  $sg \land \operatorname{DTR2} : word \land \operatorname{DTR2} : \operatorname{CAT} : v \land \operatorname{DTR2} :$  PHON : talks  $\land \operatorname{DTR2} :$  AGR :  $Y \land \operatorname{DTR2} : \operatorname{AGR} : sg \land \operatorname{DTR1} : Z_1 \land \operatorname{DTR2} : Z_2).$
- 22 word( $Z_2, 1, 2$ )  $\leftarrow X = (phrase \land CAT : s \land DTR1 : word \land DTR1 : CAT : (P 21, 4)$   $n \land DTR1 : PHON : Clinton \land DTR1 : AGR : Y \land DTR1 : AGR : sg \land DTR2 :$   $word \land DTR2 : CAT : v \land DTR2 : PHON : talks \land DTR2 : AGR : Y \land DTR2 : AGR :$  $sg \land DTR1 : Z_1 \land DTR2 : Z_2$ ).
- 23  $\operatorname{sign}(Z_2, 1, 2) \leftarrow \operatorname{word}(Z_2, 1, 2) \& X = (phrase \land CAT : np \land DTR1 : word \land (P 20, 7)$ DTR1: CAT :  $n \land DTR1$ : PHON : Clinton  $\land DTR1$ : AGR :  $sg \land DTR2$  : word  $\land DTR2$ : CAT :  $n \land DTR2$ : PHON : talks  $\land DTR1 : Z_1 \land DTR2 : Z_2$ ).
- 24 word( $Z_2, 1, 2$ )  $\leftarrow X = (phrase \land CAT : np \land DTR1 : word \land DTR1 : CAT : (P 23,5)$  $n \land DTR1 : PHON : Clinton \land DTR1 : AGR : sg \land DTR2 : word \land DTR2 : CAT :$  $n \land DTR2 : PHON : talks \land DTR2 : AGR : pl \land DTR1 : Z_1 \land DTR2 : Z_2).$
- 25  $\operatorname{sign}(Z_2, 1, 2) \leftarrow X = (phrase \land CAT : s \land DTR1 : word \land DTR1 : CAT : (C 21, 22)$   $n \land DTR1 : PHON : Clinton \land DTR1 : AGR : Y \land DTR1 : AGR : sg \land DTR2 :$   $word \land DTR2 : CAT : v \land DTR2 : PHON : talks \land DTR2 : AGR : Y \land DTR2 : AGR :$  $sg \land DTR1 : Z_1 \land DTR2 : Z_2).$
- 26 phrase $(X, 0, 2) \leftarrow X = (phrase \land CAT : s \land DTR1 : word \land DTR1 : CAT : (C 18, 25)$   $n \land DTR1 : PHON : Clinton \land DTR1 : AGR : Y \land DTR1 : AGR : sg \land DTR2 :$   $word \land DTR2 : CAT : v \land DTR2 : PHON : talks \land DTR2 : AGR : Y \land DTR2 : AGR :$  $sg \land DTR1 : Z_1 \land DTR2 : Z_2).$
- 27  $\operatorname{sign}(X, 0, 2) \leftarrow X = (phrase \land CAT : s \land DTR1 : word \land DTR1 : CAT : (C 10, 26)$   $n \land DTR1 : PHON : Clinton \land DTR1 : AGR : Y \land DTR1 : AGR : sg \land DTR2 :$   $word \land DTR2 : CAT : v \land DTR2 : PHON : talks \land DTR2 : AGR : Y \land DTR2 : AGR :$  $sg \land DTR1 : Z_1 \land DTR2 : Z_2).$
- 28  $\operatorname{sign}(Z_2, 1, 2) \leftarrow X = (phrase \land CAT : np \land DTR1 : word \land DTR1 : CAT : (C 23, 24)$  $n \land DTR1 : PHON : Clinton \land DTR1 : AGR : sg \land DTR2 : word \land DTR2 : CAT :$  $n \land DTR2 : PHON : talks \land DTR2 : AGR : pl \land DTR1 : Z_1 \land DTR2 : Z_2).$
- 29 phrase $(X, 0, 2) \leftarrow X = (phrase \land CAT : np \land DTR1 : word \land DTR1 : CAT : (C 20,28)$   $n \land DTR1 : PHON : Clinton \land DTR1 : AGR : sg \land DTR2 : word \land DTR2 : CAT :$  $n \land DTR2 : PHON : talks \land DTR2 : AGR : pl \land DTR1 : Z_1 \land DTR2 : Z_2).$
- 30  $\operatorname{sign}(X, 0, 2) \leftarrow X = (phrase \land CAT : np \land DTR1 : word \land DTR1 : CAT : (C 10, 29)$  $n \land DTR1 : PHON : Clinton \land DTR1 : AGR : sg \land DTR2 : word \land DTR2 : CAT :$  $n \land DTR2 : PHON : talks \land DTR2 : AGR : pl \land DTR1 : Z_1 \land DTR2 : Z_2).$

Figure 4.8: Earley deduction chart, cont.



Figure 4.9: Chart graph

represented by edges connecting a node with a node possibly further on its right. Such nodes are attached with the numbers of clauses derived by the completion rule.

There are two parses of the above input sentence hidden in the Earley deduction chart of Figs. 4.7 and 4.8. In the chart graph of Fig. 4.9, these two parses are represented by the upper and lower half of the symmetric graph. From each of the two final completed clauses, 27 and 30, a proof tree representing a parse can be reconstructed using the algorithm of Def. 4.6. This algorithm defines the construction of partial proof trees from completed clauses, and when applied recursively, permits the construction of proof trees from a given Earley deduction chart.

**Definition 4.6.** Let  $c_k$  be a completed clause derived from clauses  $c_i$  and  $c_j$ , let  $t_i$  and  $t_j$  be the unique partial proof trees corresponding to  $c_i$  and  $c_j$ , and define for each predicted clause  $E = (E \leftarrow F) \text{ a partial proof tree } | \quad . \text{ Then the partial proof tree } t(t_i, t_j) \text{ corresponding to } c_k \text{ is } F$ constructed s.t.

$$t(t_i, t_j) = \begin{cases} t_i \\ \oplus & if \ both \ c_i, c_j \ are \ completed \ clauses, \\ t_j \\ \\ t_i \\ \otimes & if \ one \ or \ both \ c_i, c_j \ are \ predicted \ clauses, \\ t_j \\ \\ \vdots \\ and \ \begin{array}{c} t_1 \\ \oplus \\ t_2 \end{array} = \begin{array}{c} A & A \\ B \cup C \\ t_2 \end{array} = \begin{array}{c} \vdots \\ B \cup C \\ t_2 \end{array} = \begin{array}{c} B \\ B \cup C \\ B \\ C \cup D \end{array} = \begin{array}{c} C \\ B \\ f_1 \\ B \\ f_2 \end{array} = \begin{array}{c} C \\ B \\ f_1 \\ B \\ f_2 \end{array} = \begin{array}{c} C \\ F_2$$

The proof tree for completed clause 27, corresponding to the parse  $[Clinton_N \ talks_V]_S$ , is the proof tree of Fig. 2.8 and repeated here in Fig. 4.10. The parse  $[Clinton_N \ talks_N]_{NP}$  is derived via the proof tree of Fig. 2.9, repeated here in Fig. 4.11, and can be reconstructed from completed clause 30.

### 4.9.2 Probabilistic CLGs and the Viterbi Algorithm

Turning to probabilistic CLGs, we see that because CLGs are simply instances of CLP, all techniques developed for statistical inference of probabilistic CLP apply to probabilistic CLGs



Figure 4.10: Proof tree for  $[Clinton_N \ talks_V]_S$ 



Figure 4.11: Proof tree for  $[Clinton_N \ talks_N]_{NP}$ 

without modification. The simplest way to inspect a probability distribution on parses is to list the respective proof trees and calculate their probabilities from the subtree-properties and the corresponding parameters. An imaginable probability model for the proof trees of Figs. 4.10 and 4.11 could take as properties the subtrees introduced by the clauses which are responsible for the two different readings of the input sentence, namely clauses 1, 4, 2, and 5. The respective properties  $\chi_1$ ,  $\chi_2$ ,  $\chi_3$ , and  $\chi_4$  are depicted in Figs. 4.10 and 4.11 as framed parts of the proof trees. MLE from a large natural language corpus for the parameters  $\lambda_1$ ,  $\lambda_2$ ,  $\lambda_3$  and  $\lambda_4$  corresponding to these properties would probably return a higher weight for parameters  $\lambda_1$  and  $\lambda_2$  than for  $\lambda_3$  and  $\lambda_4$ . Thus this probability model would tell the proof tree of Fig. 4.10, corresponding to the parse [*Clinton<sub>N</sub>* talks<sub>V</sub>]<sub>S</sub>, to be more probable given the input sentence *Clinton* talks than the proof tree of Fig. 4.11, corresponding to the parse [*Clinton<sub>N</sub>* talks<sub>N</sub>]<sub>NP</sub>.

However, if we are interested in the most probable parse of a sentence, listing all possible parses may be too costly in general, even if the parses just have to be extracted from a chart. Clearly, it would be nice if we could make use of the structure of the probabilistic model to guide the search for the most probable parse. The Viterbi algorithm (Viterbi (1967), Forney (1973)) for finding the most probable parse implements this idea using a dynamicprogramming approach as follows: During derivation, each derivation state must keep track of the most probable path of derivation states leading towards it. When the final derivation state is reached, the maximum probability derivation can be recovered by tracing back the stored path of most probable derivation states.

Clearly, different specifications of this algorithm depend on the chosen parsing strategy and on the underlying probability model. For example, Stolcke (1993) computes a Viterbi parse for probabilistic context-free grammars in a framework of probabilistic Earley parsing as follows: During derivation, each completed item keeps track of the most probable path of items contributing to it. The rule probabilities are propagated recursively by associating each predicted item with the probability of the rewriting rule used in the prediction, and by recording for each completed item the product of probabilities of the pair of items that contributes with maximal value to the completion. Storing at each completion step the itempair leading to the maximum, finally yields a path of most probable items from which the most probable derivation can be retrieved.

Under certain restrictions on the parsing strategy and on the probabilistic search method, the idea of the Viterbi algorithm is applicable to Earley deduction for log-linear probabilistic CLGs as follows.

Concerning the parsing strategy, let us strictly adhere to the definition of Earley deduction given above. That is, we only speak of an ambiguous derivation of a completed clause if more than one pair of clauses yields via completion the same clause with the same variable binding. That is, in this setting a numerical comparison at a completion step is done only between clause-pairs contributing via completion to the same "instantiated" clause.

Considering the probabilistic search method, Stolcke (1993)'s model of Viterbi parsing can be reconstructed if we identify the properties of the log-linear model with program clauses. If properties are allowed to be subtrees of proof trees, things are more complicated. In this setting, in order to compare numerically between alternatives, we have to incrementally build up partial proof trees and check their properties during derivation.

First, we have to define a function w to calculate the weight of a partial proof tree  $t_k$  under a log-linear probability model  $p_{\lambda}$ .

**Definition 4.7.** Let C be an Earley deduction chart for query G and program  $\mathcal{P}$ , let  $\mathcal{X}$  be the set of proof trees for G from  $\mathcal{P}$ , and let  $p_{\lambda}$  be a log-linear distribution on  $\mathcal{X}$ . Then the weight w of a partial proof tree  $t_k$  constructable for a completed clause  $c_k \in C$  is defined s.t.

$$w(t_k) = e^{\lambda \cdot \nu(t_k)}$$

Furthermore, a numerical comparison between alternatives leading to the same completion requires the partial proof trees corresponding to the alternative completions to include only completely built-up subtree-properties. This is necessary to avoid the outranking of highly weighted partial proof trees by lower weighted partial proof trees at a completion step where the highly weighted subtree-properties cannot yet be taken under consideration. For an appropriate partial ordering on trees based on an operation  $\subseteq$ , we can ensure that partial proof trees include only completely built-up properties as follows.

A partial proof tree  $t_k$  is complete for a property-vector  $\chi = (\chi_1, \ldots, \chi_n)$  iff for each  $i = 1, \ldots, n$ :  $\chi_i \subseteq t_k$  or else  $\chi_i \cap t_k = \emptyset$ .

The algorithm of Def. 4.6 can be used for a recursive comparison as follows. Note that we use the definition of variant given in Chap. 2 for the specification of an equivalence class of clauses to be compared.

For each equivalence class  $[c_k]$  of completed clauses, record the partial proof tree  $t_k^* = \underset{t_k}{\operatorname{arg max}} w(t_k)$ , where  $[c_k] = \{c \in \mathcal{C} | c \text{ is a variant of } c_k$ , and there exist clauses  $c_i$  and  $c_j$  in  $\mathcal{C}$  from which c is derivable via completion}, and  $t_k \in \{t(t_i^*, t_j^*) | t_i^* \text{ and } t_j^* \text{ are the hightest weighted complete partial proof trees corresponding to clauses <math>c_i$  and  $c_j$ }.

Clearly, given the above restrictions, this procedure will yield the most probable proof tree for a given query to a program. The possible savings in computational complexity induced by this procedure depend on the size of the subtree-properties to be worked out during the search



Figure 4.12: Type hierarchy

process. That is, small subtrees will permit an efficient pruning at nearly each completion step whereas subtrees connecting nodes over long distances may in the worst case yield no gain in efficiency at all.

#### 4.9.3 Heuristic Searching for Most Probable Parses

However, the effective applicability of the search procedure stated above strongly depends the form of the grammars under consideration. That means, for particular CLGs, it is inefficient to restrict the numerical selection only to alternative completions which lead to the same clause with the same variable binding. The storing of variable bindings in each step of an Earley deduction procedure is necessary to enable partial proofs to be reused in other partial proofs. Unfortunately, deriving a new clause with each new variable binding may introduce overhead which causes in the worst case an exponential search cost. This can be the case, e.g., for grammars which encode parses entirely via variable bindings, i.e., via  $\mathcal{L}$  -constraints, and in not via predicates, i.e.,  $\mathcal{R}(\mathcal{L})$ -atoms. The extreme ends of the spectrum of such examples can be marked, e.g., for the first case by CLGs resulting from a direct application of the compilation procedure of Götz (1995). This procedure translates HPSG descriptions into the  $\mathcal{L}$ -constraints of a CLP fragment using a single  $\mathcal{R}(\mathcal{L})$ -atom for processing. An example for the second case are definite clause grammars such as those presented in Pereira and Warren (1983) which encode each grammar symbol as a distinct CLP predicate. For cases like the first, it would be more effective if one could compare alternative completions leading to a variant of a CLP clause irrespective of the variable bindings. Unfortunately, this approach to comparing "uninstantiated" completed clauses introduces a context-dependence problem caused by incompatible variable bindings. That is, we are confronted here with a trade-off between efficiency and correctness of the search method.

Let us illustrate this context-dependence problem with a simple example. For illustration we use the program of Fig. 4.1, repeated here in Fig. 4.13, with  $\mathcal{L}$ -constraints from a language of hierarchical types. The ordering on the types is depicted in Fig. 4.12.

An Earley deduction chart for the query  $\mathbf{s}(Z) \& Z = e$  is given in Fig. 4.14.

Let the properties  $\chi_1$  to  $\chi_5$  of a probability distribution over the proof trees corresponding to the Earley deduction chart of Fig. 4.14 be defined as the framed subtrees shown in Fig. 4.15. Furthermore, let the corresponding parameter values be  $\lambda_1 = \ln 2$ ,  $\lambda_2 = \ln 3$ ,  $\lambda_3 =$ 

$$1 \ \mathbf{s}(Z) \leftarrow \mathbf{p}(Z) \ \& \ \mathbf{q}(Z)$$
$$2 \ \mathbf{p}(Z) \leftarrow Z = a.$$
$$3 \ \mathbf{p}(Z) \leftarrow Z = b.$$
$$4 \ \mathbf{p}(Z) \leftarrow Z = f.$$
$$5 \ \mathbf{q}(Z) \leftarrow Z = a.$$
$$6 \ \mathbf{q}(Z) \leftarrow Z = b.$$



7	$\leftarrow \mathbf{s}(Z) \& Z = e.$	(I)
8	$\mathbf{s}(Z) \leftarrow \mathbf{p}(Z) \ \& \ \mathbf{q}(Z) \ \& \ Z = e.$	(P 7,1)
9	$\mathbf{p}(Z) \leftarrow Z = a.$	(P 8,2)
10	$\mathbf{p}(Z) \leftarrow Z = b.$	(P 8,3)
11	$\mathbf{p}(Z) \leftarrow Z = f.$	(P 8,4)
12	$\mathbf{s}(Z) \leftarrow \mathbf{q}(Z) \ \& \ Z = a.$	(C 8,9)
13	$\mathbf{s}(Z) \leftarrow \mathbf{q}(Z) \ \& \ Z = b.$	(C 8,10)
14	$\mathbf{s}(Z) \leftarrow \mathbf{q}(Z) \ \& \ Z = f.$	(C 8,11)
15	$\mathbf{q}(Z) \leftarrow Z = a.$	(P 12,5)
16	$\mathbf{q}(Z) \leftarrow Z = b.$	(P 13,6)
17	$\mathbf{s}(Z) \leftarrow Z = a.$	(C 12,15)
18	$\mathbf{s}(Z) \leftarrow Z = b.$	(C 13,16)

Figure 4.14: Earley deduction chart

In 5,  $\lambda_4 = \ln 5$  and  $\lambda_5 = \ln 3$ . Now let us take a look at how the probability model defined by these properties and parameters guides the search for the most probable proof tree of query  $\mathbf{s}(Z) \& Z = e$  from the program of Fig. 4.13. The first decision to be made is between the completed clauses 12, 13 and 14, which differ only by their variable bindings, i.e., by their  $\mathcal{L}$ -constraints. The partial proof trees corresponding to these clauses,  $t_{12}$ ,  $t_{13}$  and  $t_{14}$ , are shown in Fig. 4.15. However,  $t_{14}$ , the highest weighted of these partial proof trees, is not included in any proof tree,  $t_{17}$  or  $t_{18}$ , corresponding to the final completion steps. That is, in this case the probabilistic search has not only missed the most probable proof tree but has led to failure! Even if we ignore the clauses contributing to this failure, namely clauses 4, 11 and 14, a problem still remains. In this case,  $w(t_{13}) > w(t_{12})$ , and the weight of proof tree  $t_{18}$  including the best partial proof tree  $t_{13}$  is  $w(t_{18}) = 9$ . However, the weight of proof tree  $t_{17}$  including the partial proof tree  $t_{12}$ , which we just have thrown away, is  $w(t_{17}) = 10$  and  $w(t_{17}) > w(t_{18})$ . Thus in this case the probabilistic search method has led us to the lower weighted proof tree.

Clearly, the first of these problems can be solved by constructing the Earley deduction chart in advance and by checking the terminal  $\mathcal{L}$ -constraint of each partial proof tree corresponding to an alternative completion against the  $\mathcal{L}$ -constraints of the final completion steps in the chart. This can be accomplished by the the following re-definition of the equivalence class  $[c_k]$ 





Figure 4.15: Partial proof trees

of completed clauses which is subject to a numerical comparison in a completion step. Note that we refer here to the definition of variable renaming given in Chap. 2.

Let  $\mathcal{C}$  be an Earley deduction chart for query G and program  $\mathcal{P}$ ,  $c_k = (A \leftarrow B_1 \& \ldots B_n \& \psi)$  be a clause in  $\mathcal{C}$ , and  $c'_k = c_k \setminus \psi$ . Then an equivalence class  $[c_k]$  of completed clauses in  $\mathcal{C}$  is defined s.t.

 $[c_k] = \{c \in \mathcal{C} | \text{ there exist clauses } c_i \text{ and } c_j \text{ in } \mathcal{C} \text{ from which clause } c = (C \leftarrow D_1 \& \dots \& D_m \& \phi) \text{ is derivable via completion, } c' = c \setminus \phi \text{ is obtained from } c'_k \text{ by simultaneously replacing each occurrence of a variable } X \text{ in } c'_k \text{ by a renamed variable } \rho(X) \text{ for all variables } X \in V(c'_k) \text{ for a renaming } \rho, \text{ and there exists a satisfiable } \mathcal{L} \text{ -constraint } \phi \& \varphi \text{ for at least one final completed clause in } \mathcal{C} \text{ with } \mathcal{L} \text{ -constraint } \varphi \}.$ 

However, the latter of the problems stated above is solvable only at the cost of reintroducing the restriction to compare only "instantiated" variants of completed clauses, i.e., variants of clauses with the same  $\mathcal{L}$ -constraints. Each search algorithm which allows "uninstantiated" variants of completed clauses to be compared, necessarily provides only a heuristic search procedure in the sense that it does not guarantee that the most probable proof tree is found. According to the definition of  $[c_k]$ , the equivalence class of completed clauses which are compared in a completion step, and the satisfaction of the completeness requirement on partial proof trees, the algorithm of Table 4.9.3 defines either an approximate heuristic or a true best-parse search algorithm.

An approach to a Viterbi-like heuristic search procedure similar to ours is used by Carroll and Briscoe (1992) for searching the parse forest produced by their probabilistic LR parser for unification-based grammars. There also the parse forest must be built up completely before unpacking to ensure that the search algorithm does pursue successful derivations.

On the whole, the decision between a possibly more efficient, but only approximate heuristic procedure and a possibly inefficient, but optimal Viterbi algorithm has to be made with respect to particular classes of CLGs in mind. Furthermore, if a heuristic search procedure is used, an alternative to completing the chart in advance is to using a backtracking procedure in connection with an incremental computation of clauses and corresponding best partial proof trees.

<b>Input</b> Log-linear model $p_{\lambda}$ on set $\mathcal{X}$ of proof trees for goal $G$ from program $\mathcal{P}$ ,
weight function $w$ , tree-constructor function $t$ , Earley deduction algorithm,
choice of equivalence class of completed clauses.
${f Output}$ Best proof tree $t_k^*$ for $G$ from ${\cal P}$ .
Procedure
Until no clauses can be added
Compute clauses by Earley deduction algorithm,
If $c_k$ is a completed clause,
Then $w^* := 0$ , compute $[c_k]$ ,
If $[c_k] = [c_l]$ for some $l < k$ ,
Then $t_k^* := t_l^*$ ,
Else for each $c \in [c_k]$ ,
For each $c_i, c_j$ which derive c via completion,
Compute the best proof tree $t_i^*$ for $c_i$ ,
Compute the best proof tree $t_j^*$ for $c_j$ ,
$t_k := t(t_i^*, t_j^*),$
${\rm If}w(t_k)>w^*,$
Then $w^* := w(t_k), t_k^* := t_k,$
Else delete $c$ ,
Return $t_k^*$ .



### 4.10 Summary and Discussion

In this chapter we presented a probabilistic model for CLP and a novel method for statistical inference about the parameters of such models from incomplete data. We discussed the problems of previous approaches which applied Baum's estimation technique for stochastic context-free models to estimation of stochastic constraint-based models. We showed with a counterexample that this incomplete-data estimation method does not generally yield the desired maximum likelihood values when applied to constraint-based systems. To overcome the inherent context-dependence problem of such systems, we introduced a powerful log-linear probability model for CLP. Furthermore, we presented a new algorithm to infer the parameters of log-linear models, and also the properties of such parametric models, from incomplete training data. We showed monotonicity and convergence of the algorithm to the desired maximum likelihood estimates and applied it experimentally to estimation of a CLG on a small scale. Furthermore, we discussed various methods for approximate computation of the formulae involved in the inference task, and presented methods which use the structure of the probabilitic model to guide the search for the most probable analysis. To this end we presented an approximate heuristic search algorithm based on dynamic programming techniques. Depending on the class of grammars under consideration, this algorithm can provide a considerable efficiency gain in searching for the most probable analysis.

In comparison to the work on quantitative CLP presented above, the advantages of probabilistic CLP are clearly the possibility to use automatic techniques for statistical inference for parameter estimation and property selection. Rather, our incomplete-data inference algorithm is general enough to be applicable to log-linear probability distributions in general, and thus is useful in other incomplete-data settings as well. In this chapter the algorithm has especially been shown to be useful for probabilistic context-sensitive NLP models. In contrast to related approaches such as that of Magerman (1994), Ratnaparkhi (1998) or Goodman (1998), which require fully annotated corpora for estimation, our statistical inference algorithm provides general means for automatic and reusable training of arbitrary probabilistic constraint-based grammars from unannotated corpora. Furthermore, our approach is the first one since the introduction of log-linear models into the discussion of probabilistic parsing by Abney (1996) which evaluates experimentally the usefulness of general log-linear models on CLGs.

## Chapter 5

# Conclusion

In this final chapter we present a short summary of the work of this thesis. We compare the advantages and shortcomings of the two presented approaches to quantitative and probablistic CLP relative to each other and relative to other approaches. Not surprisingly, the presented work is not definitive but raises several questions which could not be answered in the course of this thesis. These questions will be dealt with when we discuss future continuations of the presented work.

### 5.1 Summary

In this thesis, we have presented new mathematical and algorithmic techniques for quantitative and statistical inference in constraint-based NLP. We have chosen the general concepts of CLP as the formal framework to deal with constraint-based NLP, yielding CLGs as instances of CLP. Aiming at a general solution of the problem of structural ambiguity in CLGs, we have presented two independent approaches to weighted CLGs.

The first approach, called quantitative CLP, is situated in a clear logical framework, and presents a sound and complete system of quantitative inference for definite clauses with subjective weights attached to them. This approach permits to specify weights in arbitrary ways, e.g., as subjective probabilities, user-defined preference values, or degrees of grammaticality, and to use search techniques such as alpha-beta pruning for finding the maximally weighted proof tree for a given set of queries efficiently. Related previous work either focussed solely on formal semantics of quantitative logic programs without specific applications in mind, or presented only informal attachments of weights to grammar components for the aim of weight-based pruning in natural language parsing. Our approach is the first one to combine weight-based parsing for constraint-based systems with a rigid formal semantics for such quantitative inference systems. The second approach, called probabilistic CLP, addresses the problem of structural ambiguity resolution by a completely different form of weighted CLGs. Here a log-linear probability model is presented which defines a probability distribution over the proof trees of a constraint logic program on the basis of weights assigned to arbitrary properties of these trees. The possibility to define arbitrary features of proof trees as such properties and to estimate appropriate weights for them permits the probabilistic modeling of arbitrary context-dependencies. In this thesis we firstly evaluate empirically the applicability and feasibility of estimation of general log-linear models on CLGs. In contrast to previous approaches which were restricted to estimation from annotated data for specialized probabilistic parsing models we present an algorithm to estimate the parameters and to induce the properties of log-linear models from incomplete, unanalyzed data. The new algorithm has the same computational complexity as related complete-data inference algorithms for log-linear models. Furthermore, we address the problem of computational intractability of large summations in the inference task by discussing various techniques to approximately solve this task and present an approximate heuristic search algorithm for CLGs.

### 5.2 Future Work

As shown in Sect. 4.7, the empirical evaluation of estimating log-linear models on CLGs showed promising results both for training and evaluation on a small scale. Clearly, the main task of future work is a thorough investigation of the performance of the presented general algorithms on *larger scales* of real-world NLP applications. In larger experiments issues which were addressed so far only theoretically shall be evaluated in practice. Such issues are the empirical evaluation of *property selection*, the evaluation of various *approximation methods* in parameter estimation, or the empirical testing of the performance of non-heuristic versus *heuristic search techniques* in terms of linguistic results.

New issues which shall be addressed in larger experiments are the use of *dynamic pro*gramming techniques not only for searching for best parses but also for efficient calculation of expectations in the estimation process. Similar to the heuristic Viterbi algorithm presented for best parse search the application of dynamic programming to computing expectations will be possible only in a heuristic way. Clearly, the question to be addressed is how such heuristic estimation procedures perform in terms of linguistic evaluations.

Another issue that will become important for larger data sets is the use of reference distributions as simpler and easier to estimate *back-off* models. A reasonable choice of a reference distribution for our task is, e.g., a model defining a probability distribution on lexical-semantic head-head relations such as verb-noun pairs (see, e.g., Rooth, Riezler, Prescher, Carroll, and Beil (1999)). Such a clustering model does not require complex parsing models or costly annotated corpora, but can be estimated easily from large corpora of verb-noun pairs. Furthermore,
such a class-based model will also provide a smooth default distribution and thus help to solve the sparse data problem.

A further task of future work will be the investigation of possible applications of log-linear models and incomplete-data estimation to NLP applications different from parsing.

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