

UNIVERSITÀ DEGLI STUDI  
ROMA TRE  
Dipartimento di Filosofia



Tesi di dottorato in  
FILOSOFIA E TEORIA DELLE  
SCIENZE UMANE

UNIVERSITÉ PARIS DIDEROT  
– PARIS 7  
U.F.R. d'Informatique –  
Laboratoire PPS



Thèse de doctorat en  
SCIENCES MATHÉMATIQUES  
DE PARIS CENTRE –  
SPECIALITÉ INFORMATIQUE

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# Differential nets, experiments and reduction

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20th June 2013



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# Introduction

## Linear Logic

Starting from semantical investigations about  $\lambda$ -calculus, Girard introduced in 1986 ([Gir87]) Linear Logic (LL), a refinement of intuitionistic and classical logic allowing a fine analysis of the use of resources during the cut-elimination (i.e. execution via Curry-Howard isomorphism) process of proofs (i.e. programs via Curry-Howard correspondence: by means of the introduction of the new connectives ! and ?, LL gives a logical *status* to duplication and erasure operations (corresponding to structural rules of intuitionistic and classical logic). One of the novelties of LL is its representation of proof by means of some particular graphs, the nets, giving a more geometrical account of the cut-elimination/execution process. Of course, not all the elements of this particular class of graphs, the proof-structures, are nets, i.e. correspond to proofs of LL's sequent calculus, but there exists a correctness criterion ([DR89], among others) characterizing all (and only in certain frameworks of LL) the proof-structures which are nets. So, proof-structures become some interesting objects in themselves from a computational point of view, in virtue of their geometrical aspect. Actually, proof structures still display some sequentialized aspects because of the presence of boxes, which define erasable or duplicable resources during the cut-elimination process. A box, indeed, must contain a proof structure that satisfies some constraints: all its conclusions, called auxiliary port, except one, called principle port, must be conclusion of a ?-link. This provides a deeply inductive character to proof structures.

**Revisiting the syntax of nets.** A first contribution of my thesis is revisiting the syntax of LL nets. Inspired by the presentation of nets developed in [dCT12] according to the formalism of interaction nets introduced by Lafont in [Laf95], I defined a syntax for the multiplicative and exponential framework of LL proof structures, in which there is not an explicit constructor for boxes; so, it's possible to recover boxes by means of some functions (the arrows) in a "purely geometrical" way under suitable conditions: these arrows associate with every !-link its auxiliary doors and the cuts contained in it at depth 0. Such a geometrical approach allows to define proof structures in a non

inductive way, bringing out the Girard’s original inner meaning of Linear Logic ([Gir87]). Differently from [dCT12], my definition allows to define also the cut-elimination procedure for proof structures: the exponential reduction step involves the composition of arrows. In these proof structures I defined also a correctness criterion (connection and acyclicity) in a completely non inductive (that is “purely geometric”) way. The syntax introduced in this thesis turns out to be compatible also with the cut-elimination procedure for differential nets of  $\text{DiLL}_0$  (see below), but not those of  $\text{DiLL}$ , where a box might contain a sum of nets.

**Differential nets and Taylor expansion.** In [Ehr05] Ehrhard defines a denotational semantics for  $\lambda$ -calculus and  $\text{LL}$  proof nets: the finiteness spaces, in which types (formulas) are interpreted by vectorial spaces and  $\lambda$ -terms ( $\text{LL}$  nets) by infinitely differentiable functions defined as power series (i.e. Taylor expansions) on these spaces. Differentiation can be internalized at a syntactical level thanks to an extension of the  $\lambda$ -calculus with differential operators: the differential  $\lambda$ -calculus, introduced by Ehrhard and Regnier in [ER03]. The authors have then extended the differential operators to  $\text{LL}$  nets, obtaining the differential nets ( $\text{DiLL}$ , [ER06b]), where the differential constructors assume an interesting form: they correspond to “symmetrizing” the exponential connectives. This means that the rules related to the two modalities  $!$  and  $?$  are perfectly symmetric, apart from the promotion rule. The differential versions of  $\lambda$ -calculus and  $\text{LL}$  allow a finer analysis of the use of resources during the computation process. We call  $\Lambda^{\text{res}}$  (resp.  $\text{DiLL}_0$ ) the fragment of differential  $\lambda$ -calculus (resp.  $\text{DiLL}$ ) characterized by the fact that only linear applications are possible (resp. by the fact that there are no boxes). Linearity, that is the absence of a promotion rule, entails that every term of  $\Lambda^{\text{res}}$  (resp. every net of  $\text{DiLL}_0$ ) is trivially strongly normalizable (their sizes strictly decrease under reduction). Furthermore, linearity also causes that a term in  $\Lambda^{\text{res}}$  (resp. a  $\text{DiLL}_0$  net) reduces to a sum of terms in  $\Lambda^{\text{res}}$  (resp. a sum of  $\text{DiLL}_0$  nets), because a linear resource (it can be used exactly once), required from several parts, determines a plurality of possible choices. Hence, for every term in  $\Lambda^{\text{res}}$  (resp.  $\text{DiLL}_0$  net)  $t$ , its normal form  $\text{NF}(t)$  exists and it is a finite linear combination of  $\lambda$ -terms in  $\Lambda^{\text{res}}$  (resp.  $\text{DiLL}_0$  nets).  $\Lambda^{\text{res}}$  (resp.  $\text{DiLL}_0$ ) can be seen as an analysis tool for  $\lambda$ -calculus (resp.  $\text{LL}$  nets), thanks to the Taylor expansion  $(\ )^*$ , which associates with every  $\lambda$ -term (resp.  $\text{LL}$  net) a (potentially infinite) sum of terms in  $\Lambda^{\text{res}}$  (resp.  $\text{DiLL}_0$  nets). In [ER08] it is proved that, given a every ordinary  $\lambda$ -term  $M$ , one can sum up all the normal forms of the resource  $\lambda$ -terms in  $M^*$ . Thus one obtains the normal form  $\text{NF}(M^*)$  of  $M^*$ , a (in general) infinite linear combination of terms in  $\Lambda^{\text{res}}$  with relational coefficients. In [ER06a] it is showed that  $\text{NF}(M^*)$  is the Taylor expansion of the Böhm tree  $\text{BT}(M)$  of  $M$

(the notion of Taylor expansion is naturally extended to Böhm trees), that is

$$\text{NF}(M^*) = (\text{BT}(M))^* \quad (1)$$

In other words, the Taylor expansion commutes with normalization, where normalizing an ordinary  $\lambda$ -term means here computing its Böhm tree. The Böhm tree of a ordinary lambda term can be seen as the normal form of the head linear reduction which is the call-by-name reduction implemented by (a version of) the Krivine's abstract machine ([Kri07]).

Taking advantage of a separation theorem for differential nets ([MP07]), I have demonstrated that the Taylor expansion (without taking into account the coefficients) commutes with the cut-elimination process. This means that, for every LL net  $\pi$ , one has  $\text{NF}(\pi^*) = (\text{NF}(\pi))^*$  (the analogous of equation (1) for LL nets). At the same time Mazza and Pagani have shown two distinct  $\text{DiLL}_0$  nets  $\rho$  and  $\rho'$  in the Taylor expansion  $\pi^*$  of a LL net  $\pi$  whose respective normal forms  $\text{NF}(\rho)$  and  $\text{NF}(\rho')$  are not disjoint. This example shows the difference of the case of  $\lambda$ -calculus with respect to LL: indeed, a crucial passage in the proof of the equation (1) in [ER08, ER06a] (for the  $\lambda$ -calculus) consists in defining a coherence relation among the  $\lambda$ -terms with sources such that:

- for every ordinary  $\lambda$ -term  $M$ , all the elements of  $M^*$  are coherent two by two among them
- if  $t$  and  $t'$  are coherent  $\lambda$ -terms with sources, then  $\text{NF}(t)$  and  $\text{NF}(t')$  are disjoint.

The Mazza and Pagani's example shows that it is impossible to define such a relation on LL nets.

**Differential nets and experiments of relational semantics.** In my thesis I tried to understand precisely and rigorously the strict relationship between differential nets without boxes (i.e. resources  $\lambda$ -terms) and experiments of LL nets. An experiment (notion introduced in [Gir87] and studied in detail in [Tor00, Tor03]) is a function which permits to associate with every LL net  $\pi$  a point of the interpretation of  $\pi$  in the relational model, the interpretation of  $\pi$  being the set of points resulting from all the possible experiments of  $\pi$ . Experiments, hence, act as a bridge between syntax and semantics. Among all the points of relational semantics of a LL net  $\pi$ , some of them are "more important": the injective points, that is those in which every their atom occurs exactly twice. If  $\pi$  is without cuts, the injective points are the results of experiments that associates with every axiom a different element of the web. From the injective points it is possible to reconstruct every other point by substitution. Furthermore two injective points of the interpretation of  $\pi$  can uniquely differ for the atom names, showing the "same structure".

Then, we say that the two injective points are equivalent and that the one can be transformed in the other by a suitable substitution of atoms with atoms. Given a net  $\pi$  of LL, we denote by  $\llbracket \pi \rrbracket$  the subset of the relational interpretation of  $\pi$ , formed by the injective points quotiented modulo the injective substitutions. In collaboration with Tortora de Falco and Pellissier, I have demonstrated that the Taylor expansion of a cut-free  $\eta$ -expanded LL net  $\pi$  coincides with  $\llbracket \pi \rrbracket$ . In other words, a differential net in the Taylor expansion of a cut-free  $\eta$ -expanded LL net  $\pi$  is a canonical representative of an equivalence class of injective points of the relational interpretation of  $\pi$ . It remains to be investigated the meaning of a differential net in the Taylor expansion of a LL net  $\pi$  with cuts. In this case the difficulty is that it can reduce itself in several differential nets without cuts.

In collaboration with Tortora de Falco and Pellissier I have characterized the relations of the relational model corresponding to interpretations of some acyclic and connected LL net. In fact, if two cut-free  $\eta$ -expanded LL nets  $\pi_1$  and  $\pi_2$ , acyclic and connected have the same 2-point in the Taylor expansion (i.e. the the differential net obtained recursively taking for every box two copies of its contents), then  $\pi_1 = \pi_2$ . In other words, a cut-free,  $\eta$ -expanded, acyclic and connected LL net is completely characterized by the 2-point in its Taylor expansion. This result simplifies the proof of injectivity of relational semantics with respect to LL (see [dCT12]) in the acyclic and connected case. Moreover thanks to this result it is possible to define an algorithm that, given a relation of the relational model, takes its 2-point  $\alpha$  (if it exists) and tries to recover a cut-free,  $\eta$ -expanded, acyclic and connected LL net. If this procedure ends successfully, then one has found the only cut-free and  $\eta$ -expanded LL net that has  $\alpha$  in its relational interpretation; otherwise, no acyclic and connected LL net has  $\alpha$  in its interpretation. This result of surjectivity is based on the fundamental hypothesis of connection, pointing out the importance of this notion.

## Call-by-value lambda-calculus

In the ordinary (also called “call-by-name”)  $\lambda$ -calculus, the prototype of any functional programming language, the values are either variables or abstractions ( $\lambda$ -terms of the shape  $\lambda xM$ ). So, the  $\lambda$ -terms are either values or applications ( $\lambda$ -terms of the shape  $MN$ ). The “call-by-value”  $\lambda$ -calculus is the version of  $\lambda$ -calculus allowing to reduce only the  $\beta_v$ -redexes, i.e.  $\beta$ -redexes of the shape  $(\lambda xM)V$  where  $V$  is a value. The call-by-value  $\lambda$ -calculus was introduced by Plotkin in '70 ([Pl075]) in order to give a version of  $\lambda$ -calculus closer to the real implementation of functional programming languages. The relationship between call-by-value  $\lambda$ -calculus and Linear Logic was widely studied for the first time by Maraist, Wadler *et al.* in [MOTW95] in '90.

Recently in [Ehr12] Ehrhard introduced a version (called  $\Lambda_{CBV}$ ) of the call-by-value  $\lambda$ -calculus such that values and terms are disjoint sets defined



by mutual induction: a value is either a variable or an abstraction  $\lambda xM$  where  $M$  is a term, a term is either an application  $(M)N$  where  $M$  and  $N$  are terms, or a “promoted” value  $V^!$  where  $V$  is a value. This distinction can be explained from the Linear Logic point of view: in [Ehr12] Ehrhard presented a general notion of denotational model for  $\Lambda_{\text{CBV}}$  corresponding to the translation  $(\ )^b$  defined as “boring” by Girard ([Gir87]) of the intuitionistic logic into LL, whereby  $(A \Rightarrow B)^b = !A^b \multimap !B^b$  (thus in the untyped case, the intuitionistic isomorphism  $o \simeq (o \Rightarrow o)$  becomes  $o \simeq (!o \multimap !o)$ ). In my thesis I studied the relationship between  $\Lambda_{\text{CBV}}$  and LL from a syntactical point of view (already implicit in [Ehr12]). I defined the translation of terms and values of  $\Lambda_{\text{CBV}}$  into LL nets: the idea is that a “promoted” value corresponds to a box in the LL nets, therefore a  $\beta_v$ -redex  $(\lambda xM)^!V^!$  corresponds to a cut between the box representing  $(\lambda xM)^!$  and a dereliction (the application is linear on the left); this allows the box representing  $V^!$  to get in the net representing  $M$  and duplicate at will. In general, one step of  $\beta_v$ -reduction in  $\Lambda_{\text{CBV}}$  corresponds to several steps of cut-elimination in LL-nets.

**Reduction and call-by-value Krivine’s machine. Trees.** In [Ehr12] Ehrhard proved that the interpretation of a term  $M$  in  $\Lambda_{\text{CBV}}$  is empty if and only if  $M$  is strongly normalizable for the  $\hat{\beta}_v$ -reduction, where the  $\hat{\beta}_v$ -reduction (or weak  $\beta_v$ -reduction) is a restriction of the  $\beta_v$ -reduction obtained by forbidding reductions under abstractions. This result is the analogous of the well-known theorem for the ordinary (i.e. call-by-name)  $\lambda$ -calculus whereby a term is head normalizable if and only if its interpretation in the Engler’s denotational model is not empty. In my thesis I developed a survey about  $\hat{\beta}_v$ -reduction, in order to see to what extent the  $\hat{\beta}_v$ -reduction can be considered in  $\Lambda_{\text{CBV}}$  as an analogue of the head reduction in ordinary  $\lambda$ -calculus. A first difference is obvious: in the ordinary  $\lambda$ -calculus the head redex of any term, if any, is unique, whereas a term in  $\Lambda_{\text{CBV}}$  might have several  $\hat{\beta}_v$ -redexes (in LL-nets they correspond to cuts at depth 0), but these  $\hat{\beta}_v$ -redexes are not overlapping, hence the  $\hat{\beta}_v$ -reduction is strongly confluent. Therefore, one can define a parallel  $\hat{\beta}_v$ -reduction reducing in one step all the  $\hat{\beta}_v$ -redexes: if a term  $M$  is  $\hat{\beta}_v$ -normalizable, then the parallel  $\hat{\beta}_v$ -reduction reduces  $M$  to its  $\hat{\beta}_v$ -normal form. So, the fact of having several  $\hat{\beta}_v$ -redexes is not a substantial difference with respect to the head reduction of ordinary  $\lambda$ -calculus.

The structure of a term  $M$  can be represented by a binary tree, called the applicative tree of  $M$ : it breaks up the applications in  $M$  until to “promoted” values which are subterms of  $M$  (they are the leaves of the applicative tree of  $M$ ). So, the  $\hat{\beta}_v$ -redexes of any term are characterizable as the nodes whose left (resp. right) child is a “promoted” abstraction (resp. “promoted” value). The notions of applicative tree and parallel  $\hat{\beta}_v$ -reduction suggest the definition of a tree-like structure which is similar to a Böhm tree for the “call-by-value”

$\lambda$ -calculus. Nowadays for the “call-by-value”  $\lambda$ -calculus there does not yet exist a notion of Böhm tree (see for example the recent [NGP12]).

In my thesis I also defined an abstract machine for  $\Lambda_{\text{CBV}}$  similar to the Krivine’s abstract machine for the ordinary (i.e. call-by-name)  $\lambda$ -calculus define in [Kri07, DR04]. The abstract machines play an important role in implementing programming languages because on the one hand they are “sufficiently abstract” to relate easily to the notion of reduction of  $\lambda$ -calculus, on the other hand they are closer to executions of a real machine, by imposing among other things a precise reduction strategy. I introduced two call-by-value Krivine’s machines  $K^l$  and  $K^r$ : I showed that the  $K^l$  (resp.  $K^r$ ) machine with an input term  $M$  will search for the leftmost (resp. rightmost)  $\hat{\beta}_v$ -redex in the applicative tree of  $M$  and then reduce it. By the good proprieties of the  $\hat{\beta}_v$ -reduction, if a closed term  $M$  is  $\hat{\beta}_v$ -normalizable then its  $\hat{\beta}_v$ -normal form computed by  $K^l$  and  $K^r$ ; actually this result holds more generally for any “random” call-by-value Krivine’s machine at each execution step chooses whether to apply the left or right reduction strategy.

**Translations,  $\sigma$ -equivalence and  $\sigma$ -reduction.** There exists two continuation passing style (CPS) translations of  $\Lambda_{\text{CBV}}$  into the ordinary (i.e. call-by-name)  $\lambda$ -calculus:  $(\ )^l$  (already defined in [Plo75, Sel01]) and  $(\ )^r$  (introduced in my thesis), whose only difference is in the translation of application, more precisely in the choice of putting the function (in case of  $(\ )^l$ ) or the argument (in case of  $(\ )^r$ ). I showed that, modulo these CPS translations, the  $\beta_v$ -reduction corresponds to the  $\beta$ -reduction of ordinary  $\lambda$ -calculus. The following result is more interesting: the call-by-value Krivine’s machine  $K^l$  (resp.  $K^r$ ) is simulated by the call-by-name Krivine’s machine modulo the CPS translation  $(\ )^l$  (resp.  $(\ )^r$ ).

In the ordinary  $\lambda$ -calculus,  $\sigma$ -equivalence ([Reg92, Reg94]) equates terms differing only in their sequential structure but behaving the same. The  $\sigma$ -equivalence can be characterized by encoding  $\lambda$ -terms into LL nets by means of the Girard’s “call-by-name” translation  $(A \Rightarrow B) \rightsquigarrow (!A \multimap B)$ : two  $\lambda$ -terms are  $\sigma$ -equivalents if and only if their translations into LL nets are the same. I proved an analogous result for  $\Lambda_{\text{CBV}}$  by means of the “boring” translation of the intuitionistic arrow  $(A \Rightarrow B) \rightsquigarrow (!A \multimap !B)$  into LL. The  $\sigma_v$ -equivalence relation thus obtained on terms and values in  $\Lambda_{\text{CBV}}$  is not included in  $\beta_v$ -equivalence, differently from what happens in ordinary  $\lambda$ -calculus, where  $\sigma$ -equivalence is included in  $\beta$ -equivalence.

Even more surprisingly, I showed that it is possible to give an orientation to two of the three rules generating  $\sigma_v$ -equivalence, in such a way to get a “completion” of  $\beta_v$ -reduction: the add of the  $\sigma$ -reduction rules allows to simulate the Accattoli and Paolini’s call-by-value  $\lambda$ -calculus with explicit substitutions ( $\lambda_{\text{vsub}}$ , introduced in [AP12]). One of the novelties of  $\lambda_{\text{vsub}}$  is that it allows to characterize the solvable terms by means of internal (i.e.

call-by-value) reduction rules. Thanks to the simulation, it is reasonable to expect that the solvability is characterizable in  $\Lambda_{\text{CBV}}$  by means of internal (i.e. call-by-value) reduction rules without using explicit substitutions.



Part I  
Linear Logic



# Chapter 1

## A non inductive syntax

This section is devoted to present in full details the syntactical object for which we prove our main result: proof-structure (definition 44). We adopt the interaction nets point of view (see for example [Laf95, ER06b, Pag09, Tra11, dCT12]) and pass through intermediate objects: cell-bases (definition 1), pre-pre-proof-structures (definition 12), pre-proof-structures (definition 35). Our approach, definitions and notations are those of [dCT12] (in particular, our syntactical objects are untyped as in [LT06, PT10, dCPT11]) up to some differences that will be explained in the following. Essentially the principal novelties with respect to the syntax of [dCT12] are:

- our framework can represent DiLL-proof-structures, which are the differential generalization (where boxes and duals of  $?$ -links are allowed, see for example [ER06b, MP07, Pag09, Tra11]) of the MELL-proof-structures (the multiplicative-exponential framework of Linear Logic, see for example [Gir87, DR95, Tor03, dCPT11, dCT12]);
- our objects are not necessarily cut-free; moreover it is possible to define the cut-elimination in two frameworks of our syntax, the DiLL<sub>0</sub>-proof-structures (the DiLL-proof-structures without boxes) and the MELL-proof-structures; this fact answers positively to the difficulties raised in [dCT12] about the definition of cut-elimination on untyped proof-structures;
- our definition of proof-structure is completely non-inductive, so a proof-structure is precisely a labeled hyper-graph; the boxes are computed by starting from its principal door and by using only some “geometrical informations” in this hyper-graph; our geometrical point of view is strengthened by our choice of untyped syntactical objects.

**Notation.** We set  $\mathcal{T} = \{1, \perp, \otimes, \wp, !, ?\}$  whose elements are the *connectives of the multiplicative and exponential framework of Linear Logic*. We say that  $1, \perp, \otimes, \wp$  (resp.  $!, ?$ ) are the *multiplicative* (resp. *exponential*) connectives, and  $1, \perp$  are the *units*.

## 1.1 Cells and ports

In the following definition of cell-base, we introduce cells and ports. This definition differs from that one in [dCT12] only because in our cell-base there is not the function  $\#$ : this means that the word “linear” used in [dCT12] makes no sense in our syntax.

**Definition 1** (Module-base, (pseudo-)cell-base). *A module-base is a 5-tuple  $\mathbb{C} = (\mathbf{t}, \mathcal{P}, \mathbf{C}, \mathbf{P}^{\text{pri}}, \mathbf{P}^{\text{left}})$  such that:*

- $\mathbf{t}$  is a function such that  $\text{dom}(\mathbf{t})$  is a finite set and  $\text{codom}(\mathbf{t}) = \mathcal{T} \cup \{ax\}$ ; we set  $\mathcal{C}(\mathbb{C}) = \text{dom}(\mathbf{t})$  whose elements are the cells of  $\mathbb{C}$ ; for every  $l \in \mathcal{C}(\mathbb{C})$ ,  $\mathbf{t}(l)$  is the label of  $l$ ; for every  $t, t' \in \mathcal{T}$ , we set  $\mathcal{C}^t(\mathbb{C}) = \{l \in \mathcal{C}(\mathbb{C}) \mid \mathbf{t}(l) = t\}$  (whose elements are the  $t$ -cells of  $\mathbb{C}$ ),  $\mathcal{C}^{t,t'}(\mathbb{C}) = \mathcal{C}^t(\mathbb{C}) \cup \mathcal{C}^{t'}(\mathbb{C})$  and  $\mathcal{C}_2^{\otimes, \mathfrak{A}}(\mathbb{C}) = \{l \in \mathcal{C}^{\otimes, \mathfrak{A}}(\mathbb{C}) \mid \mathbf{a}_{\mathbb{C}}(l) = 2\}$ ;
- $\mathcal{P}$  is a finite set whose elements are the ports of  $\mathbb{C}$ ; we set  $\mathcal{P}(\mathbb{C}) = \mathcal{P}$ ;
- $\mathbf{C} : \mathcal{P}(\mathbb{C}) \rightarrow \mathcal{C}(\mathbb{C})$  is a surjection such that for every  $l \in \mathcal{C}(\mathbb{C})$ ,
  - if  $\mathbf{t}(l) \in \{1, \perp, ax\}$  then  $\text{card}(\{p \in \mathcal{P}(\mathbb{C}) \mid \mathbf{C}(p) = l\}) = 1$ ,
  - if  $\mathbf{t}(l) \in \{\otimes, \mathfrak{A}\}$  then  $1 \leq \text{card}(\{p \in \mathcal{P}(\mathbb{C}) \mid \mathbf{C}(p) = l\}) \leq 3$ ;

for every  $l \in \mathcal{C}(\mathbb{C})$ , we set  $\mathcal{P}_l(\mathbb{C}) = \{p \in \mathcal{P}(\mathbb{C}) \mid \mathbf{C}(p) = l\}$  whose elements are the ports of  $l$ ;

- $\mathbf{P}^{\text{pri}} : \mathcal{C}(\mathbb{C}) \rightarrow \mathcal{P}(\mathbb{C})$  is a function such that  $\mathbf{C} \circ \mathbf{P}^{\text{pri}} = \text{id}_{\mathcal{C}(\mathbb{C})}$ ; for every  $l \in \mathcal{C}(\mathbb{C})$ ,  $\mathbf{P}^{\text{pri}}(l)$  is the principal port (or conclusion) of  $l$ , moreover we set  $\mathcal{P}_l^{\text{aux}}(\mathbb{C}) = \mathcal{P}_l(\mathbb{C}) \setminus \{\mathbf{P}^{\text{pri}}(l)\}$  whose elements are the auxiliary ports (or premises) of  $l$ , and  $\mathbf{a}_{\mathbb{C}}(l) = \text{card}(\mathcal{P}_l^{\text{aux}}(\mathbb{C}))$  which is the arity of  $l$ ; we set  $\mathcal{P}^{\text{pri}}(\mathbb{C}) = \text{im}(\mathbf{P}^{\text{pri}})$  whose elements are the principal ports of  $\mathbb{C}$ , and  $\mathcal{P}^{\text{aux}}(\mathbb{C}) = \bigcup_{l \in \mathcal{C}(\mathbb{C})} \mathcal{P}_l^{\text{aux}}(\mathbb{C})$  whose elements are the auxiliary ports of  $\mathbb{C}$ ;
- $\mathbf{P}^{\text{left}} : \mathcal{C}_2^{\otimes, \mathfrak{A}}(\mathbb{C}) \rightarrow \mathcal{P}^{\text{aux}}(\mathbb{C})$  is a function such that, for every  $l \in \mathcal{C}_2^{\otimes, \mathfrak{A}}(\mathbb{C})$ , one has  $\mathbf{P}^{\text{left}}(l) \in \mathcal{P}_l^{\text{aux}}(\mathbb{C})$ .

A pseudo-cell-base is a module-base such that  $\mathcal{C}^{\otimes, \mathfrak{A}}(\mathbb{C}) = \mathcal{C}_2^{\otimes, \mathfrak{A}}(\mathbb{C})$ .

A cell-base is a pseudo-cell-base such that  $\mathcal{C}^{ax}(\mathbb{C}) = \emptyset$ .

We denote by **ModuleBases** (resp. **PseudoCells**; **Cells**) the set of module-bases (resp. pseudo-cell-bases; cell-bases).

Intuitively, a module-base corresponds to a set of “links with their premises and conclusions” in the standard theory of linear logic proof-nets (see for example [Gir87, DR95, Tor03, Pag09, dCPT11]). More precisely, cells correspond to links, the principal port of a cell corresponds to the conclusion of a link and an auxiliary port of a cell corresponds to a premise of a link.



Note that our presentation reformulates the linear logic “*nouvelle syntaxe*” of [Reg92, DR95] (where the ?-links have any arity) in the style of (differential) interaction nets (see [Laf95, ER06b]).

**Notation.** Let  $\mathbb{C} = (\mathfrak{t}, \mathcal{P}, \mathbb{C}, \mathcal{P}^{\text{pri}}, \mathcal{P}^{\text{left}})$  be a module-base. We set  $\mathfrak{t}_{\mathbb{C}} = \mathfrak{t}$ ,  $\mathbb{C}_{\mathbb{C}} = \mathbb{C}$ ,  $\mathcal{P}_{\mathbb{C}}^{\text{pri}} = \mathcal{P}^{\text{pri}}$ ,  $\mathcal{P}_{\mathbb{C}}^{\text{left}} = \mathcal{P}^{\text{left}}$ . We recall the notations  $\mathcal{P}(\mathbb{C}) = \mathcal{P}$  and  $\mathcal{C}(\mathbb{C}) = \text{dom}(\mathfrak{t}_{\mathbb{C}})$ .

For  $\mathbb{C} \in \mathbf{ModuleBases}$ , the function  $\mathcal{P}_{\mathbb{C}}^{\text{pri}}$  allows to distinguish the principal port from the auxiliary ones of any cell of  $\mathbb{C}$ . As expected, for the binary  $\otimes$ - and  $\wp$ -cells of  $\mathbb{C}$ , the function  $\mathcal{P}_{\mathbb{C}}^{\text{left}}$  allows to distinguish the left auxiliary port from the right one, whereas for the other kinds of cells a similar function is not defined because their auxiliary ports (if any) are not ordered.

Typically a cell  $l$  in a module-base  $\mathbb{C}$  is graphically depicted as a triangle with its label  $\mathfrak{t}_{\mathbb{C}}(l)$  inside, the principal port being on a vertex and the auxiliary ones on the opposed side (in such a way that when the principal port is downwards the left auxiliary ports of a binary  $\otimes$ - or  $\wp$ -cell is placed on the left).

**Remark 2.** Let  $\mathbb{C} \in \mathbf{ModuleBases}$ .

The functions  $\mathcal{P}_{\mathbb{C}}^{\text{pri}}$  and  $\mathcal{P}_{\mathbb{C}}^{\text{left}}$  induce the functions:

- $\mathcal{P}_{\mathbb{C}}^{\text{aux}} : \mathcal{C}(\mathbb{C}) \rightarrow \mathcal{P}(\mathcal{P}^{\text{aux}}(\mathbb{C}))$  defined by  $\mathcal{P}_{\mathbb{C}}^{\text{aux}}(l) = \mathcal{P}_l^{\text{aux}}$  for every  $l \in \mathcal{C}(\mathbb{C})$ ; thus  $\text{im}(\mathcal{P}_{\mathbb{C}}^{\text{aux}}) = \mathcal{P}^{\text{aux}}(\mathbb{C})$ ;
- $\mathcal{P}_{\mathbb{C}}^{\text{right}} : \mathcal{C}_2^{\otimes, \wp}(\mathbb{C}) \rightarrow \mathcal{P}^{\text{aux}}(\mathbb{C})$  defined by  $\{\mathcal{P}_{\mathbb{C}}^{\text{right}}(l)\} = \mathcal{P}_l^{\text{aux}} \setminus \{\mathcal{P}_{\mathbb{C}}^{\text{left}}(l)\}$  for every  $l \in \mathcal{C}_2^{\otimes, \wp}(\mathbb{C})$ ; note that  $\mathcal{P}_{\mathbb{C}}^{\text{right}}$  is well-defined since the binary  $\otimes$ - and  $\wp$ -cells have exactly two auxiliary ports.

Notice that  $\mathcal{P}(\mathbb{C}) = \mathcal{P}^{\text{pri}}(\mathbb{C}) \uplus \mathcal{P}^{\text{aux}}(\mathbb{C})$  and  $\mathcal{P}_l(\mathbb{C}) = \mathcal{P}_l^{\text{pri}}(\mathbb{C}) \uplus \mathcal{P}_l^{\text{aux}}(\mathbb{C})$  for every  $l \in \mathcal{C}(\mathbb{C})$ .

Furthermore, if  $\mathbb{C} \in \mathbf{PseudoCells}$ , then  $\text{dom}(\mathcal{P}_{\mathbb{C}}^{\text{left}}) = \mathcal{C}_2^{\otimes, \wp}(\mathbb{C}) = \text{dom}(\mathcal{P}_{\mathbb{C}}^{\text{right}})$ .

Among the cell-bases, there is the *empty cell-base*  $\mathbb{C}$  defined by  $\mathcal{P}(\mathbb{C}) = \emptyset$  and  $\mathfrak{t}_{\mathbb{C}}$ ,  $\mathbb{C}_{\mathbb{C}}$ ,  $\mathcal{P}_{\mathbb{C}}^{\text{pri}}$  and  $\mathcal{P}_{\mathbb{C}}^{\text{left}}$  are empty functions.

The following notion will be used in definitions 4 and 6

**Definition 3** (Completeness). *Let  $\mathbb{C} \in \mathbf{ModuleBases}$  and let  $Q \subseteq \mathcal{P}(\mathbb{C})$ :  $Q$  is  $\mathbb{C}$ -complete when, for every  $l \in \mathcal{C}(\mathbb{C})$ , if  $\mathcal{P}_{\mathbb{C}}^{\text{pri}}(l) \in Q$ , then  $\mathcal{P}_l^{\text{aux}}(\mathbb{C}) \subseteq Q$ .*

The following definitions 4, 6 and 8 formalize some intuitive notions of:

- erasure of some cells and ports in a module-base;
- submodule-base of a module-base;
- disjoint union of module-bases.

**Definition 4** (Erasure of cells and ports). *Let  $\mathbb{C} \in \mathbf{ModuleBases}$  and  $n \in \mathbb{N}$ .*

*Let  $l_1, \dots, l_n \in \mathcal{C}(\mathbb{C})$ , let  $Q \subseteq \mathcal{P}(\mathbb{C})$  be  $\mathbb{C}$ -complete and let  $L_Q = \{l \in \mathcal{C}(\mathbb{C}) \mid \mathcal{P}_l(\mathbb{C}) \subseteq Q\}$ . The erasure of  $l_1, \dots, l_n$  and  $Q$  in  $\mathbb{C}$  is  $\mathbb{C}' = (\mathfrak{t}', \mathcal{P}', \mathbb{C}', \mathcal{P}'^{\text{pri}}, \mathcal{P}'^{\text{left}})$  where:*

- $\mathfrak{t}' = \mathfrak{t}_{\mathbb{C}} \setminus \mathcal{C}(\mathbb{C}) \setminus (L_Q \cup \{l_1, \dots, l_n\})$ ;
- $\mathcal{P}' = \mathcal{P}(\mathbb{C}) \setminus (Q \cup \bigcup_{i=1}^n \mathcal{P}_{l_i}(\mathbb{C}))$ ;
- $\mathbb{C}' = \mathbb{C}_{\mathbb{C}} \setminus \mathcal{P}'$ ;
- $\mathcal{P}'^{\text{pri}} = \mathcal{P}_{\mathbb{C}}^{\text{pri}} \setminus \mathcal{C}(\mathbb{C}) \setminus (L_Q \cup \{l_1, \dots, l_n\})$ ;
- $\mathcal{P}'^{\text{left}} = \mathcal{P}_{\mathbb{C}}^{\text{left}} \setminus \mathcal{C}_2^{\otimes, \mathfrak{A}}(\mathbb{C}) \setminus (L_Q \cup \{l_1, \dots, l_n\})$ .

We say then that “ $\mathbb{C}'$  is obtained from  $\mathbb{C}$  by erasing  $l_1, \dots, l_n$  and  $Q$ ”.

**Remark 5.** For every  $\mathbb{C} \in \mathbf{ModuleBases}$  (resp.  $\mathbb{C} \in \mathbf{PseudoCells}$ ;  $\mathbb{C} \in \mathbf{Cells}$ ) and  $l_1, \dots, l_n \in \mathcal{C}(\mathbb{C})$ , if  $\mathbb{C}'$  is obtained from  $\mathbb{C}$  by erasing  $l_1, \dots, l_n$ , then  $\mathbb{C}' \in \mathbf{ModulesBases}$  (resp.  $\mathbb{C}' \in \mathbf{PseudoCells}$ ;  $\mathbb{C}' \in \mathbf{Cells}$ ); moreover, if  $Q \subseteq \mathcal{P}(\mathbb{C})$  is  $\mathbb{C}$ -complete and  $\mathbb{C}'$  is obtained from  $\mathbb{C}$  by erasing  $l_1, \dots, l_n$  and  $Q$ , then  $\mathbb{C}' \in \mathbf{ModulesBases}$ .

**Definition 6** (Submodule-base). *Let  $\mathbb{C} \in \mathbf{ModuleBases}$ . Let  $Q \subseteq \mathcal{P}(\mathbb{C})$  be  $\mathbb{C}$ -complete and such that, for every  $l \in \mathcal{C}(\mathbb{C})$  and  $p \in Q$ , if  $p \in \mathcal{P}_l^{\text{aux}}(\mathbb{C})$  then  $\mathcal{P}_l(\mathbb{C}) \subseteq Q$ ;<sup>1</sup> let  $L_Q = \{l \in \mathcal{C}(\mathbb{C}) \mid \mathcal{P}_l(\mathbb{C}) \subseteq Q\}$ . The submodule-base of  $\mathbb{C}$  generated by  $Q$  is  $\text{module}_{\mathbb{C}}(Q) = (\mathfrak{t}', \mathcal{P}', \mathbb{C}', \mathcal{P}'^{\text{pri}}, \mathcal{P}'^{\text{left}})$  where:*

- $\mathfrak{t}' = \mathfrak{t}_{\mathbb{C}} \setminus L_Q$ ;
- $\mathcal{P}' = Q$ ;
- $\mathbb{C}' = \mathbb{C}_{\mathbb{C}} \setminus Q$ ;
- $\mathcal{P}'^{\text{pri}} = \mathcal{P}_{\mathbb{C}}^{\text{pri}} \setminus L_Q$ ;
- $\mathcal{P}'^{\text{left}} = \mathcal{P}_{\mathbb{C}}^{\text{left}} \setminus \mathcal{C}_2^{\otimes, \mathfrak{A}}(\mathbb{C}) \cap L_Q$ .

**Remark 7.** Let  $\mathbb{C} \in \mathbf{ModuleBases}$  (resp.  $\mathbb{C} \in \mathbf{PseudoCells}$ ;  $\mathbb{C} \in \mathbf{Cells}$ ). Let  $Q \subseteq \mathcal{P}(\mathbb{C})$  be  $\mathbb{C}$ -complete and such that, for every  $l \in \mathcal{C}(\mathbb{C})$  and  $p \in Q$ , if  $p \in \mathcal{P}_l^{\text{aux}}(\mathbb{C})$  then  $\mathcal{P}_l(\mathbb{C}) \subseteq Q$ . Then  $\text{module}_{\mathbb{C}}(Q) \in \mathbf{Modules}$  (resp.  $\text{module}_{\mathbb{C}}(Q) \in \mathbf{PseudoCells}$ ;  $\text{module}_{\mathbb{C}}(Q) \in \mathbf{Cells}$ ).

<sup>1</sup>According to definition 3, this entails that, for every  $l \in \mathcal{C}(\mathbb{C})$ , either  $\mathcal{P}_l(\mathbb{C}) \cap Q = \emptyset$  or  $\mathcal{P}_l(\mathbb{C}) \subseteq Q$ .

**Definition 8** (Disjoint union of module-bases). *Let  $\mathbb{C}$  and  $\mathbb{C}' \in \mathbf{ModuleBases}$ :  $\mathbb{C}$  and  $\mathbb{C}'$  are disjoint if  $\mathcal{C}(\mathbb{C}) \cap \mathcal{C}(\mathbb{C}') = \emptyset$  and  $\mathcal{P}(\mathbb{C}) \cap \mathcal{P}(\mathbb{C}') = \emptyset$ .*

*Let  $n \in \mathbb{N}$  and let  $\mathbb{C}_1, \dots, \mathbb{C}_n \in \mathbf{ModuleBases}$  be pairwise disjoint: the disjoint union of  $\mathbb{C}_1, \dots, \mathbb{C}_n$  is*

$$\biguplus_{i=1}^n \mathbb{C}_i = \left( \bigcup_{i=1}^n \mathfrak{t}_{\mathbb{C}_i}, \bigcup_{i=1}^n \mathcal{P}(\mathbb{C}_i), \bigcup_{i=1}^n \mathbb{C}_{\mathbb{C}_i}, \bigcup_{i=1}^n \mathcal{P}_{\mathbb{C}_i}^{\text{pri}}, \bigcup_{i=1}^n \mathcal{P}_{\mathbb{C}_i}^{\text{left}} \right).$$

*If  $n = 2$ , the disjoint union of  $\mathbb{C}_1$  and  $\mathbb{C}_2$  is denoted by  $\mathbb{C}_1 \uplus \mathbb{C}_2$ .*

**Remark 9.** For every  $n \in \mathbb{N}$ , if  $\mathbb{C}_1, \dots, \mathbb{C}_n \in \mathbf{ModuleBases}$  (resp.  $\mathbb{C}_1, \dots, \mathbb{C}_n \in \mathbf{PseudoCells}$ ;  $\mathbb{C}_1, \dots, \mathbb{C}_n \in \mathbf{Cells}$ ) are pairwise disjoint and  $\mathbb{C} = \biguplus_{i=1}^n \mathbb{C}_i$ , then  $\mathbb{C} \in \mathbf{ModuleBases}$  (resp.  $\mathbb{C} \in \mathbf{PseudoCells}$ ;  $\mathbb{C} \in \mathbf{Cells}$ ).

We introduce the notion of “identity” (or better said isomorphism) between two module-bases. The idea is that two module-bases are isomorphic iff they are identical up to the names of their cells and ports (in particular, they have the same graphical representation).

**Definition 10** (Isomorphism on module-bases). *Let  $\mathbb{C}, \mathbb{C}' \in \mathbf{ModuleBases}$ .*

*An isomorphism from  $\mathbb{C}$  to  $\mathbb{C}'$  is a pair  $\varphi = (\varphi_{\mathcal{C}}, \varphi_{\mathcal{P}})$  of bijections  $\varphi_{\mathcal{C}} : \mathcal{C}(\mathbb{C}) \rightarrow \mathcal{C}(\mathbb{C}')$  and  $\varphi_{\mathcal{P}} : \mathcal{P}(\mathbb{C}) \rightarrow \mathcal{P}(\mathbb{C}')$  such that the following diagrams commute:*

$$\begin{array}{ccc} \mathcal{C}(\mathbb{C}) & \xrightarrow{\mathcal{P}_{\mathbb{C}}^{\text{pri}}} & \mathcal{P}(\mathbb{C}) & \xrightarrow{\mathbb{C}_{\mathbb{C}}} & \mathcal{C}(\mathbb{C}) & \xrightarrow{\mathfrak{t}_{\mathbb{C}}} & \mathcal{T} \\ \varphi_{\mathcal{C}} \downarrow & & \varphi_{\mathcal{P}} \downarrow & & \varphi_{\mathcal{C}} \downarrow & \nearrow \mathfrak{t}_{\mathbb{C}'} & \\ \mathcal{C}(\mathbb{C}') & \xrightarrow{\mathcal{P}_{\mathbb{C}'}^{\text{pri}}} & \mathcal{P}(\mathbb{C}') & \xrightarrow{\mathbb{C}_{\mathbb{C}'}} & \mathcal{C}(\mathbb{C}') & & \end{array} \quad \begin{array}{ccc} \mathcal{C}_2^{\otimes, \mathfrak{R}}(\mathbb{C}) & \xrightarrow{\mathcal{P}_{\mathbb{C}}^{\text{left}}} & \mathcal{P}(\mathbb{C}) \\ \varphi_{\mathcal{C}} \downarrow & & \downarrow \varphi_{\mathcal{P}} \\ \mathcal{C}_2^{\otimes, \mathfrak{R}}(\mathbb{C}') & \xrightarrow{\mathcal{P}_{\mathbb{C}'}^{\text{left}}} & \mathcal{P}(\mathbb{C}') \end{array}$$

*We write then  $\varphi : \mathbb{C} \simeq \mathbb{C}'$ .*

*If there exists an isomorphism from  $\mathbb{C}$  to  $\mathbb{C}$ , then we say that  $\mathbb{C}$  and  $\mathbb{C}'$  are isomorphic and we write  $\mathbb{C} \simeq \mathbb{C}'$ .*

**Remark 11.** Let  $\mathbb{C}, \mathbb{C}' \in \mathbf{ModuleBases}$ . If  $\varphi$  is an isomorphism from  $\mathbb{C}$  to  $\mathbb{C}'$  then:

1.  $\text{im}(\varphi_{\mathcal{P}} \upharpoonright_{\mathcal{P}^{\text{aux}}(\mathbb{C})}) = \mathcal{P}^{\text{aux}}(\mathbb{C}')$  and  $\text{im}(\varphi_{\mathcal{P}} \upharpoonright_{\mathcal{P}_l^{\text{aux}}(\mathbb{C})}) = \mathcal{P}_{\varphi_{\mathcal{C}}(l)}^{\text{aux}}(\mathbb{C}')$  (in particular,  $\mathfrak{a}_{\mathbb{C}}(l) = \mathfrak{a}_{\mathbb{C}'}(\varphi_{\mathcal{C}}(l))$ ) for every  $l \in \mathcal{C}(\mathbb{C})$ ;
2. if  $\mathbb{C} \in \mathbf{PseudoCells}$  (resp.  $\mathbb{C} \in \mathbf{Cells}$ ) then  $\mathbb{C}' \in \mathbf{PseudoCells}$  (resp.  $\mathbb{C}' \in \mathbf{Cells}$ ).

## 1.2 Pre-pre-proof-structures

A pre-pre-proof-structure (see also the analogous definition in [dCT12]) is morally a (hyper-)graph consisting of a cell-base, isolated ports (not belonging to any cell of the cell-base), wires connecting the ports of its cells and the isolated ones, and arrows to add some informations.

**Definition 12** (Module, (pseudo-)pre-pre-proof-structure). *A module is a 6-tuple  $\Phi = (\mathbb{C}, \mathcal{I}, \mathcal{D}, \mathcal{W}, \text{auxd}, \text{bc})$  where:*

- $\mathbb{C} \in \mathbf{ModuleBases}$  is the module-base of  $\Phi$ ; we set  $\mathcal{C}(\Phi) = \mathcal{C}(\mathbb{C})$  whose elements are the cells of  $\Phi$ ;
- $\mathcal{I}$  and  $\mathcal{D}$  are finite sets (whose elements are respectively the isolated ports of  $\Phi$  and the deadlocks of  $\Phi$ ), satisfying  $\mathcal{I} \cap \mathcal{P}(\mathbb{C}) = \emptyset$ ,  $\mathcal{D} \cap \mathcal{P}(\mathbb{C}) = \emptyset$  and  $\mathcal{I} \cap \mathcal{D} = \emptyset$ ; we set  $\mathcal{P}(\Phi) = \mathcal{P}(\mathbb{C}) \cup \mathcal{I} \cup \mathcal{D}$  whose elements are the ports of  $\Phi$ ;  $\Phi$  is deadlock-free if  $\mathcal{D}(\Phi) = \emptyset$ ;
- $\mathcal{W} \subseteq \mathcal{P}_2(\mathcal{P}(\Phi) \setminus \mathcal{D})$  such that:
  1. for every  $w, w' \in \mathcal{W}$ , if  $w \cap w' \neq \emptyset$  then  $w = w'$ ,
  2.  $\mathcal{P}^{\text{aux}}(\mathbb{C}) \cup \mathcal{I} \subseteq \bigcup \mathcal{W}$ ,

the elements of  $\mathcal{W}$  are the wires of  $\Phi$ ; we set  $\mathcal{Cuts}(\Phi) = \{\{p, q\} \in \mathcal{W} \mid p, q \in \mathcal{P}^{\text{pri}}(\mathbb{C})\}$  whose elements are the cuts of  $\Phi$ ; any  $p \in \bigcup \mathcal{Cuts}(\Phi)$  is a cut port of  $\Phi$ ;  $\Phi$  is cut-free if  $\mathcal{Cuts}(\Phi) = \emptyset$ ;

- $\text{auxd}$  is a partial function from  $\mathcal{C}^!(\mathbb{C})$  to  $\mathcal{P}(\mathcal{P}^{\text{aux}}(\mathbb{C}))$  such that for every  $l \in \mathcal{C}(\mathbb{C})$ , if  $\text{auxd}$  is defined in  $l$  then:
  - $\text{a}_{\mathbb{C}}(l) = 1$ ,
  - if  $p \in \text{auxd}(l)$ , then  $p \in \mathcal{P}_l^{\text{aux}}(\mathbb{C})$  for some  $l' \in \mathcal{C}^2(\mathbb{C})$ ; we say that  $p$  is an auxiliary door of  $l$ ;

we set  $\mathcal{C}^{\text{prom}}(\Phi) = \text{dom}(\text{auxd})$  (resp.  $\text{Auxdoors}(\Phi) = \bigcup \text{im}(\text{auxd})$ ) whose elements are the promotion cells (resp. auxiliary doors) of  $\Phi$ ; if  $l \in \mathcal{C}^{\text{prom}}(\Phi)$  and  $p$  is the premise of  $l$ , we set  $\text{doors}_{\Phi}(l) = \{p\} \cup \text{auxd}(l)$  and we say that  $p$  is the principal door of  $l$  in  $\Phi$  and any  $q \in \text{auxd}(l)$  is an auxiliary door of  $l$  in  $\Phi$ ;

- $\text{bc}$  is a function from  $\mathcal{C}^{\text{prom}}(\Phi)$  to  $\mathcal{P}(\bigcup \mathcal{Cuts}(\Phi) \cup \mathcal{D})$  such that:
  - if  $\text{bc}(l) \cap \text{bc}(l') \neq \emptyset$  for some  $l, l' \in \mathcal{C}^{\text{prom}}(\Phi)$ , then  $l = l'$ ,<sup>2</sup>
  - if  $\{p, q\} \in \mathcal{Cuts}(\Phi)$  and  $p \in \text{bc}(l) \cap \bigcup \mathcal{Cuts}(\Phi)$  for some  $l \in \mathcal{C}^{\text{prom}}(\Phi)$ , then  $q \in \text{bc}(l)$ ,<sup>3</sup>

we set  $\text{cutports}_{\Phi}(l) = \text{bc}(l) \cap \bigcup \mathcal{Cuts}(\Phi)$  and  $\text{deadlocks}_{\Phi}(l) = \text{bc}(l) \cap \mathcal{D}$  for every  $l \in \mathcal{C}^{\text{prom}}(\Phi)$ .

<sup>2</sup>This conditions means that for every cut or deadlock, there exists at most one !-cell pointing to it.

<sup>3</sup>This conditions means that, for every cut  $w$  of  $\Phi$ , the function  $\text{bc}$  either points to both ports of  $w$  or does not point to any port of  $w$ . Therefore, we are entitled to talk about a cut associated with a promotion cell by the function  $\text{bc}$ .

A pseudo-structure is a module  $\Phi = (\mathbb{C}, \mathcal{I}, \mathcal{D}, \mathcal{W}, \text{auxd}, \text{bc})$  such that  $\mathbb{C}$  is a pseudo-cell-base,  $\{\{p, p'\} \in \mathcal{W} \mid \exists l, l' \in \mathcal{C}^{ax}(\mathbb{C}) : p = \mathbb{P}_{\mathbb{C}}^{\text{pri}}(l) \text{ and } p' = \mathbb{P}_{\mathbb{C}}^{\text{pri}}(l')\} = \emptyset^4$  and

3. for every  $w \in \mathcal{W}$ , if  $w \cap \mathcal{I} \neq \emptyset$  then  $w \cap \mathcal{P}^{\text{pri}}(\mathbb{C}) = \emptyset$ ;

we say then that  $\mathbb{C}$  is the pseudo-cell-base of  $\Phi$ .

A pre-pre-proof-structure (or ppps for short) is a pseudo-structure  $\Phi$  such that the pseudo-cell-base  $\mathbb{C}$  of  $\Phi$  is a cell-base; we say then that  $\mathbb{C}$  is the cell-base of  $\Phi$ .

We denote by **Modules** (resp. **PseudoPPPS**; **PPPS**) the set of modules (resp. pseudo-structures; pre-pre-proof-structures).

In a module, an isolated port is depicted as a dot, a wire  $\{p, q\}$  is graphically depicted as a line connecting the ports  $p$  and  $q$ , a deadlock is graphically depicted as a circle. If  $l$  is a promotion cell then its label is depicted as  $!p$ , furthermore the fact that an auxiliary port  $q$  of a  $?$ -cell is an auxiliary door of  $l$  is represented graphically by a dotted arrow from  $l$  to  $q$ ; likewise, if  $q$  is a deadlock or cut port in  $\text{bc}(l)$ , this is represented graphically by a dotted arrow from  $l$  to  $q$  or to the cut  $w$  such that  $q \in w$ .

A promotion cell of a ppps  $\Phi$  has to be seen as a “candidate for a box”, i.e. a cell which is the starting point to attempt to compute the box (a particular sub-graph of  $\Phi$ ) associated with it (in general, it is not always possible, see definition 38).

Our definition of pre-pre-proof-structure differs from that one in [dCT12] by the following points:

- in our pre-pre-proof-structures, cuts (wires connecting the principal ports of two different cells) are allowed;
- in order to be closed under cut-elimination, in our definition of ppps we add the set  $\mathcal{D}$  of deadlocks; a deadlock has to be seen as a sort of degenerate cut (morally, it is an axiom whose conclusions are connected by a cut, but our syntax cannot express that explicitly);
- in order to handle differential nets with our syntax, the  $!$ -cells’ arity does not need to be 1; furthermore, not all unary  $!$ -cells are “candidates for a box”;
- with respect to the definition in [dCT12], in our ppps we add the “arrow” functions  $\text{auxd}$  and  $\text{bc}$  which associate with every promotion cell  $l$  respectively the set of its auxiliary doors and the set of cut ports and deadlocks of depth<sup>5</sup> 0 in the “box-candidate” associated with  $l$ ;

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<sup>4</sup>This means that in a pseudo-structure there is no cut connecting the principal ports of two  $ax$ -cells.

<sup>5</sup>See definition 51 and proposition 54 for the notion of depth.

a promotion cell  $l$  might have no auxiliary doors (resp. no cuts nor deadlocks) associated with it, this is the case when  $\text{auxd}(l) = \emptyset$  (resp.  $\text{bc}(l) = \emptyset$ ).

**Definition 13** (Free port, axiom, arrow). *Let  $\Phi = (\mathbb{C}, \mathcal{I}, \mathcal{D}, \mathcal{W}, \text{auxd}, \text{bc}) \in \mathbf{Modules}$ . We set:*

- $\mathbb{C}(\Phi) = \mathbb{C}$ ,  $\mathcal{I}(\Phi) = \mathcal{I}$ ,  $\mathcal{W}(\Phi) = \mathcal{W}$ ,  $\mathcal{D}(\Phi) = \mathcal{D}$ ,  $\text{auxd}_\Phi = \text{auxd}$ ,  $\text{bc}_\Phi = \text{bc}$ ;  $\mathcal{C}^t(\Phi) = \mathcal{C}^t(\mathbb{C}(\Phi))$  and  $\mathcal{C}^{t,t'}(\Phi) = \mathcal{C}^{t,t'}(\mathbb{C}(\Phi))$  for every  $t, t' \in \mathcal{T}$ ;  $\mathcal{C}_2^{\otimes, \mathfrak{A}}(\Phi) = \mathcal{C}_2^{\otimes, \mathfrak{A}}(\mathbb{C}(\Phi))$ ;  $\mathcal{P}^{\text{pri}}(\Phi) = \mathcal{P}^{\text{pri}}(\mathbb{C}(\Phi))$ ,  $\mathcal{P}^{\text{aux}}(\Phi) = \mathcal{P}^{\text{aux}}(\mathbb{C}(\Phi))$ ;  $\mathcal{P}_l^{\text{aux}}(\Phi) = \mathcal{P}_l^{\text{aux}}(\mathbb{C}(\Phi))$  and  $\mathcal{P}_l(\Phi) = \mathcal{P}_l(\mathbb{C}(\Phi))$  for every  $l \in \mathcal{C}(\Phi)$ ;  $\mathbb{P}_\Phi^{\text{pri}} = \mathbb{P}_{\mathbb{C}(\Phi)}^{\text{pri}}$ ,  $\mathbb{P}_\Phi^{\text{left}} = \mathbb{P}_{\mathbb{C}(\Phi)}^{\text{left}}$ ,  $\mathbb{P}_\Phi^{\text{aux}} = \mathbb{P}_{\mathbb{C}(\Phi)}^{\text{aux}}$ ;  $\mathfrak{t}_\Phi = \mathfrak{t}_{\mathbb{C}(\Phi)}$ ,  $\mathbb{C}_\Phi = \mathbb{C}_{\mathbb{C}(\Phi)}$ ,  $\mathfrak{a}_\Phi = \mathfrak{a}_{\mathbb{C}(\Phi)}$ ;
- $\mathcal{P}^{\text{free}}(\Phi) = \mathcal{I}(\Phi) \cup (\mathcal{P}^{\text{pri}}(\Phi) \setminus \bigcup \mathcal{W}(\Phi))$  whose elements are the free ports (or conclusions) of  $\Phi$ ;  $\mathcal{C}^{\text{term}}(\Phi) = \{l \in \mathcal{C}(\Phi) \mid \mathbb{P}_\Phi^{\text{pri}}(l) \in \mathcal{P}^{\text{free}}(\Phi)\}$  whose elements are the terminal cells of  $\Phi$ ;
- $\mathcal{Ax}(\Phi) = \{\{p, q\} \in \mathcal{W}(\Phi) \mid p, q \notin \mathcal{P}^{\text{pri}}(\Phi)\}$  whose elements are the axioms of  $\Phi$ ; any  $p \in \bigcup \mathcal{Ax}(\Phi)$  is an axiom port of  $\Phi$ ;  $\mathcal{Ax}^{\text{term}}(\Phi) = \{w \in \mathcal{Ax}(\Phi) \mid \exists p \in w : p \in \mathcal{I}(\Phi)\}$  (resp.  $\mathcal{Ax}^{\text{isol}}(\Phi) = \{w \in \mathcal{Ax}(\Phi) \mid \forall p \in w : p \in \mathcal{I}(\Phi)\}$ ) whose elements are the terminal (resp. isolated) axioms of  $\Phi$ ;
- $\mathcal{Arrows}(\Phi) = \{\{p, q\} \in \mathcal{P}_2(\mathcal{P}(\Phi)) \mid \exists l \in \mathcal{C}^{\text{prom}}(\Phi) : p \in \mathcal{P}_l^{\text{aux}}(\Phi), q \in \text{auxd}_\Phi(l) \cup \text{bc}_\Phi(l)\}$ , whose elements are the arrows of  $\Phi$ ;
- $\mathcal{Cuts}_0(\Phi) = \mathcal{Cuts}(\Phi) \setminus \mathcal{P}_2(\bigcup \text{im}(\text{bc}_\Phi))$  (whose elements are the cuts at depth 0 of  $\Phi$ ) and  $\mathcal{D}_0(\Phi) = \mathcal{D}(\Phi) \setminus \bigcup \text{im}(\text{bc}_\Phi)$  (whose elements are the deadlocks at depth 0 of  $\Phi$ ).

For a module  $\Phi$ ,  $p$  is an isolated port of  $\Phi$  when  $p$  is a port of some axiom and a conclusion of  $\Phi$ . The meaning of the conditions on the set of wires in definition 12 is the following:

- condition 1 implies that three ports cannot be connected by two wires,
- condition 2 entails that auxiliary ports can never be conclusions of a ppps,
- condition 3 (only for pseudo-structures) implies that when the principal port of some cell is connected to another port this is necessarily a port of some cell, hence “hanging” wires (i.e. connecting a principal port and an isolated one) are not allowed in pseudo-structures.

Intuitively, a module  $\Phi$  can be seen as:

- a finite undirected graph whose labeled nodes are the cells, deadlocks and free ports of  $\Phi$ , and whose edges are the wires of  $\Phi$  and the arrows connecting each promotion cell of  $\Phi$  with its auxiliary doors, its cut ports and its deadlocks;
- a finite undirected hyper-graph whose nodes are the ports and deadlocks of  $\Phi$ , whose labeled hyper-edges are the cells (connecting all its ports) of  $\Phi$ , and whose edges are the wires of  $\Phi$  and arrows connecting each promotion cell of  $\Phi$  with its auxiliary doors, its cut ports and its deadlocks.

**Remark 14.** Let  $\Phi \in \mathbf{Modules}$ .  $\mathcal{P}(\Phi) = \mathcal{P}^{\text{pri}}(\Phi) \uplus \mathcal{P}^{\text{aux}}(\Phi) \uplus \mathcal{I}(\Phi) \uplus \mathcal{D}(\Phi)$  and  $\mathcal{P}(\Phi) \setminus \bigcup \mathcal{W}(\Phi) = \{\mathcal{P}_{\Phi}^{\text{pri}}(l) \mid l \in \mathcal{C}^{\text{term}}(\Phi)\} \uplus \mathcal{D}(\Phi)$ , thus any port of  $\Phi$  is a conclusion of  $\Phi$  iff it is either the principal port of a terminal cell of  $\Phi$  or an axiom port of  $\Phi$ . In particular,  $\mathcal{I}(\Phi) \subseteq \bigcup \mathcal{Ax}(\Phi)$ .

Among the ppps, the *empty ppps*  $\Phi$  is defined by:

- $\mathbb{C}(\Phi)$

The following notion defines how to transform two different *ax*-cells of a pseudo-structure in an axiom: it will be used to associate with every point of  $D^{<\omega}$  a DiLL<sub>0</sub>-proof-structure (see definitions 35, 69 and 83).

**Definition 15** (Connecting pairs of *ax*-cells). *Let  $\Phi \in \mathbf{PseudoPPPS}$ .*

*Let  $l_1, l_2 \in \mathcal{C}^{\text{ax}}(\Phi)$  with  $l_1 \neq l_2$ . We say that  $\Phi'$  is obtained from  $\Phi$  by connecting  $l_1$  and  $l_2$  if  $\Phi' = (\mathbb{C}', \mathcal{I}', \mathcal{D}', \mathcal{W}', \text{auxd}', \text{bc}')$  where  $\mathcal{D}' = \mathcal{D}(\Phi)$ ,  $\text{auxd}' = \text{auxd}_{\Phi}$ ,  $\text{bc}' = \text{bc}_{\Phi}$ ,  $\mathbb{C}'$  is obtained by  $\mathbb{C}(\Phi)$  by erasing  $l_1$  and  $l_2$ , and furthermore:*

- *if  $l_1$  and  $l_2$  are not terminal cells of  $\Phi$  and  $p_1$  and  $p_2$  are the auxiliary port of  $\Phi$  such that  $\{\mathcal{P}_{\Phi}^{\text{pri}}(l_i), p_i\} \in \mathcal{W}(\Phi)$  for  $i \in \{1, 2\}$ , then  $\mathcal{W}' = (\mathcal{W}(\Phi) \setminus \{\{\mathcal{P}_{\Phi}^{\text{pri}}(l_1), p_1\}, \{\mathcal{P}_{\Phi}^{\text{pri}}(l_2), p_2\}\}) \cup \{\{p_1, p_2\}\}$  and  $\mathcal{I}' = \mathcal{I}(\Phi)$ ;*
- *if  $l_1$  is a terminal cell of  $\Phi$  and  $l_2$  is not and  $p_2$  is the auxiliary port of  $\Phi$  such that  $\{\mathcal{P}_{\Phi}^{\text{pri}}(l_2), p_2\} \in \mathcal{W}(\Phi)$ , then  $\mathcal{I}' = \mathcal{I}(\Phi) \cup \{\mathcal{P}_{\Phi}^{\text{pri}}(l_1)\}$  and  $\mathcal{W}' = (\mathcal{W}(\Phi) \setminus \{\{\mathcal{P}_{\Phi}^{\text{pri}}(l_2), p_2\}\}) \cup \{\{\mathcal{P}_{\Phi}^{\text{pri}}(l_1), p_2\}\}$ ;*
- *if  $l_2$  is a terminal cell of  $\Phi$  and  $l_1$  is not and  $p_1$  is the auxiliary port of  $\Phi$  such that  $\{\mathcal{P}_{\Phi}^{\text{pri}}(l_1), p_1\} \in \mathcal{W}(\Phi)$ , then  $\mathcal{I}' = \mathcal{I}(\Phi) \cup \{\mathcal{P}_{\Phi}^{\text{pri}}(l_2)\}$  and  $\mathcal{W}' = (\mathcal{W}(\Phi) \setminus \{\{\mathcal{P}_{\Phi}^{\text{pri}}(l_1), p_1\}\}) \cup \{\{\mathcal{P}_{\Phi}^{\text{pri}}(l_2), p_1\}\}$ ;*
- *if  $l_1$  and  $l_2$  are terminal cells of  $\Phi$ , then  $\mathcal{I}' = \mathcal{I}(\Phi) \cup \{\mathcal{P}_{\Phi}^{\text{pri}}(l_1), \mathcal{P}_{\Phi}^{\text{pri}}(l_2)\}$  and  $\mathcal{W}' = \mathcal{W}(\Phi) \cup \{\{\mathcal{P}_{\Phi}^{\text{pri}}(l_2), \mathcal{P}_{\Phi}^{\text{pri}}(l_1)\}\}$ .*

*Let  $n \in \mathbb{N}$  and let  $l_1, l'_1, \dots, l_n, l'_n$  be pairwise distinct *ax*-cells of  $\Phi$ . We say that  $\Phi'$  is obtained from  $\Phi$  by connecting  $(l_1, l'_1), \dots, (l_n, l'_n)$  when:*

- *if  $n = 0$  then  $\Phi' = \Phi$ ;*

- if  $n > 0$  then  $\Phi'$  is obtained from  $\Phi''$  by connecting  $l_n$  and  $l'_n$ , where  $\Phi''$  is obtained from  $\Phi$  by connecting  $(l_1, l'_1), \dots, (l_{n-1}, l'_{n-1})$ .

**Remark 16.** For every  $\Phi \in \mathbf{PseudoPPPS}$  and pairwise distinct  $ax$ -cells  $l_1, l'_1, \dots, l_n, l'_n$  (for some  $n \in \mathbb{N}$ ), if  $\Phi'$  is obtained from  $\Phi$  by connecting  $(l_1, l'_1), \dots, (l_n, l'_n)$  then  $\Phi' \in \mathbf{PseudoPPPS}$ ; moreover, if  $\{l_1, l'_1, \dots, l_n, l'_n\} = \mathcal{C}^{ax}(\Phi)$ , then  $\Phi' \in \mathbf{PPPS}$ .

WHY?

The following notion will be used in ??.

**Definition 17** (Erasure of terminal cells, erasure of a cut at depth 0, erasure of hanging wires). *Let  $\Phi \in \mathbf{Modules}$  and let  $n \in \mathbb{N}$ .*

*Let  $l_1, \dots, l_n \in \mathcal{C}^{\text{term}}(\Phi)$  be such that, for every  $v \in \mathcal{C}^{\text{prom}}(\Phi)$  and  $1 \leq i \leq n$ , one has  $\mathcal{P}_{l_i}^{\text{aux}}(\Phi) \cap \text{auxd}_{\Phi}(v) = \emptyset$ . The erasure of  $l_1, \dots, l_n$  in  $\Phi$  is  $\Phi' = (\mathbb{C}', \mathcal{I}', \mathcal{D}', \mathcal{W}', \text{auxd}', \text{bc}')$  where:*

- $\mathbb{C}'$  is obtained from  $\mathbb{C}(\Phi)$  by erasing  $l_1, \dots, l_n$ ;
- $\mathcal{I}' = \mathcal{I}(\Phi) \cup \{p \in \bigcup_{i=1}^n \mathcal{P}_{l_i}^{\text{aux}}(\Phi) \mid \exists q \in \mathcal{P}(\Phi) \setminus \mathcal{P}^{\text{pri}}(\Phi) : \{p, q\} \in \mathcal{W}(\Phi)\}$ ;
- $\mathcal{D}' = \mathcal{D}(\Phi)$ ;
- $\mathcal{W}' = \{\{p, q\} \in \mathcal{W}(\Phi) \mid p \notin \mathcal{P}^{\text{pri}}(\Phi) \text{ or } q \notin \bigcup_{i=1}^n \mathcal{P}_{l_i}^{\text{aux}}(\Phi)\}$ ;
- $\text{auxd}' = \text{auxd}_{\Phi} \upharpoonright_{\mathcal{C}^{\text{prom}}(\Phi) \setminus \{l_1, \dots, l_n\}}$ ;
- $\text{bc}' = \text{bc}_{\Phi} \upharpoonright_{\mathcal{C}^{\text{prom}}(\Phi) \setminus \{l_1, \dots, l_n\}}$ .

*We say then that “ $\Phi'$  is obtained from  $\Phi$  by erasing  $l_1, \dots, l_n$ ”.*

*Let  $w_1, \dots, w_n \in \text{Cuts}_0(\Phi)$ . The erasure of  $w_1, \dots, w_n$  in  $\Phi$  is  $\Phi' = (\mathbb{C}', \mathcal{I}', \mathcal{D}', \mathcal{W}', \text{auxd}', \text{bc}')$  where:*

- $\mathbb{C}' = \mathbb{C}(\Phi)$ ,  $\mathcal{I}' = \mathcal{I}(\Phi)$  and  $\mathcal{D}' = \mathcal{D}(\Phi)$ ;
- $\mathcal{W}' = \mathcal{W}(\Phi) \setminus \{w_1, \dots, w_n\}$ ;
- $\text{auxd}' = \text{auxd}_{\Phi}$  and  $\text{bc}' = \text{bc}_{\Phi}$ .

*We say then that “ $\Phi'$  is obtained from  $\Phi$  by erasing  $w_1, \dots, w_n$ ”.*

*Let  $H = \{p \in \mathcal{P}^{\text{free}}(\Phi) \mid \exists q \in \mathcal{P}^{\text{pri}}(\Phi) : \{p, q\} \in \mathcal{W}(\Phi)\}$ . The erasure of the hanging wires in  $\Phi$  is  $\text{nohang}(\Phi) = (\mathbb{C}', \mathcal{I}', \mathcal{D}', \mathcal{W}', \text{auxd}', \text{bc}')$  where:*

- $\mathbb{C}' = \mathbb{C}(\Phi)$ ;
- $\mathcal{I}' = \mathcal{I}(\Phi) \setminus H$ ;
- $\mathcal{D}' = \mathcal{D}(\Phi)$ ;
- $\mathcal{W}' = \mathcal{W}(\Phi) \setminus \{w \in \mathcal{W}(\Phi) \mid \exists p \in w \cap H\}$ ;
- $\text{auxd}' = \text{auxd}_{\Phi}$  and  $\text{bc}' = \text{bc}_{\Phi}$ .



Let  $T$  be a set such that  $T \cap (\mathcal{P}^{\text{free}}(\Phi) \cap \mathcal{P}^{\text{pri}}(\Phi)) = \emptyset$  and  $\rho : \mathcal{P}^{\text{free}}(\Phi) \cap \mathcal{P}^{\text{pri}}(\Phi) \rightarrow T$  be a bijection. The add of the hanging wires in  $\Phi$  is  $\text{hang}(\Phi) = (\mathbb{C}', \mathcal{I}', \mathcal{D}', \mathcal{W}', \text{auxd}', \text{bc}')$  where:

- $\mathbb{C}' = \mathbb{C}(\Phi)$ ;
- $\mathcal{I}' = \mathcal{I}(\Phi) \cup T$ ;
- $\mathcal{D}' = \mathcal{D}(\Phi)$ ;
- $\mathcal{W}' = \mathcal{W}(\Phi) \cup \{\{q, \rho(q)\} \mid q \in \mathcal{P}^{\text{free}}(\Phi) \cap \mathcal{P}^{\text{pri}}(\Phi)\}$ ;
- $\text{auxd}' = \text{auxd}_\Phi$  and  $\text{bc}' = \text{bc}_\Phi$ .

**Remark 18.** Let  $\Phi \in \mathbf{Modules}$  (resp.  $\Phi \in \mathbf{PseudoPPPS}$ ;  $\Phi \in \mathbf{PPPS}$ ) and let  $n \in \mathbb{N}$ .

Let  $l_1, \dots, l_n \in \mathcal{C}^{\text{term}}(\Phi)$  be such that, for every  $v \in \mathcal{C}^{\text{prom}}(\Phi)$  and  $1 \leq i \leq n$ , one has  $\mathcal{P}_{l_i}^{\text{aux}}(\Phi) \cap \text{auxd}_\Phi(v) = \emptyset$ . If  $\Phi'$  is the erasure of  $l_1, \dots, l_n$  in  $\Phi$  then  $\Phi' \in \mathbf{Modules}$  (resp.  $\Phi' \in \mathbf{PseudoPPPS}$ ;  $\Phi' \in \mathbf{PPPS}$ ).

Let  $w_1, \dots, w_n \in \mathcal{Cuts}_0(\Phi)$ . If  $\Phi'$  is the erasure of  $w_1, \dots, w_n$  in  $\Phi$  then  $\Phi' \in \mathbf{Modules}$  (resp.  $\Phi' \in \mathbf{PseudoPPPS}$ ;  $\Phi' \in \mathbf{PPPS}$ ).

One has  $\text{nohang}(\Phi), \text{hang}(\Phi) \in \mathbf{Modules}$ . Furthermore, if  $\Phi \in \mathbf{PseudoPPPS}$  then  $\text{nohang}(\Phi) = \Phi$ .

Roughly speaking, the erasure of a terminal cell  $l$  in a module  $\Phi$  is the module obtained from  $\Phi$  by erasing  $l$ , its principal port, any hanging wire created by this erasure and the auxiliary ports of  $l$  which are not axiom ports in  $\Phi$ . This operation might create new isolated ports: the auxiliary ports of  $l$  which are axiom ports of  $\Phi$ . The request that no auxiliary port of  $l$  is pointed by an arrow of any promotion cell of  $\Phi$  is mandatory to make sure that the erasure of  $l$  in  $\Phi$  is a module.

The following notions will be used in definitions 20 and 22.

**Definition 19** (Completeness and erasability of a set of ports). *Let  $\Phi \in \mathbf{Modules}$ .*

*Let  $Q \subseteq \mathcal{P}(\Phi)$ :  $Q$  is  $\Phi$ -complete (resp.  $\Phi$ -erasable) if  $Q$  is  $\mathbb{C}(\Phi)$ -complete and such that, for every  $\{p, q\} \in \mathcal{W}(\Phi)$ , if  $p \in Q \setminus \mathcal{P}^{\text{pri}}(\Phi)$  (resp. if  $p \in Q$ ) then  $q \in Q$ .*

The following definitions 20, 22 and 24 formalize some intuitive notions of:

- erasure of some ports, cells and wires in a module;
- submodule of a module;
- disjoint union of modules (it will be used in definition 83).

They are generalizations to modules of the corresponding operations seen in definition 4, 6 and 8.

**Definition 20** (Erasure of ports, cells and wires). *Let  $\Phi \in \mathbf{Modules}$  and let  $Q \subseteq \mathcal{P}(\Phi)$  be  $\Phi$ -erasable.*

*The erasure of  $Q$  in  $\Phi$  is  $\Phi' = (\mathbb{C}', \mathcal{I}', \mathcal{D}', \mathcal{W}', \text{auxd}', \text{bc}')$  where:*

- $\mathbb{C}'$  is the erasure of  $Q$  in  $\Phi$ ;
- $\mathcal{I}' = \mathcal{I}(\Phi) \setminus Q$ ;
- $\mathcal{D}' = \mathcal{D}(\Phi) \setminus Q$ ;
- $\mathcal{W}' = \{w \in \mathcal{W}(\Phi) \mid \forall p \in w : p \notin Q\}$ ;<sup>6</sup>
- $\text{auxd}' : (\mathcal{C}^{\text{prom}}(\Phi) \setminus L_Q) \rightarrow \mathcal{P}(\mathcal{P}^{\text{aux}}(\Phi))$  is a function such that, for every  $l \in \mathcal{C}^{\text{prom}}(\Phi) \setminus L_Q$ , one has  $\text{auxd}'(l) = \text{auxd}_\Phi(l) \setminus Q$ ;
- $\text{bc}' : (\mathcal{C}^{\text{prom}}(\Phi) \setminus L_Q) \rightarrow \mathcal{P}(\bigcup \mathcal{Cuts}(\Phi) \cup \mathcal{D})$  is a function such that, for every  $l \in \mathcal{C}^{\text{prom}}(\Phi) \setminus L_Q$ , one has  $\text{bc}'(l) = \text{bc}_\Phi(l) \setminus Q$ .

We say then that “ $\Phi'$  is obtained from  $\Phi$  by erasing  $Q$ ”.

**Remark 21.** Let  $\Phi \in \mathbf{Modules}$ . If  $Q \subseteq \mathcal{P}(\Phi)$  is  $\Phi$ -erasable and  $\Phi'$  is the erasure of  $Q$  in  $\Phi$ , then  $\Phi' \in \mathbf{Modules}$ .

**Definition 22** (Submodule). *Let  $\Phi \in \mathbf{Modules}$ .*

*Let  $Q \subseteq \mathcal{P}(\Phi)$  be  $\Phi$ -complete, let  $L_Q = \{l \in \mathcal{C}(\Phi) \mid \mathcal{P}_l(\Phi) \subseteq Q\}$  and let  $Q' = \{p \in Q \mid \exists l \in L_Q : p \in \mathcal{P}_l(\Phi)\}$ . The submodule of  $\Phi$  generated by  $Q$  is  $\text{module}_\Phi(Q) = (\mathbb{C}', \mathcal{I}', \mathcal{D}', \mathcal{W}', \text{auxd}', \text{bc}')$  where:*

- $\mathbb{C}' = \text{module}_{\mathbb{C}(\Phi)}(Q')$ ;
- $\mathcal{I}' = Q \setminus (Q' \cup \mathcal{D}(\Phi))$ ;
- $\mathcal{D}' = \mathcal{D}(\Phi) \cap Q$ ;
- $\mathcal{W}' = \{w \in \mathcal{W}(\Phi) \mid \forall p \in w : p \in Q\}$ ;
- $\text{auxd}' : (\mathcal{C}^{\text{prom}}(\Phi) \cap L_Q) \rightarrow \mathcal{P}(\mathcal{P}^{\text{aux}}(\Phi))$  is a function such that, for every  $l \in \mathcal{C}^{\text{prom}}(\Phi) \cap L_Q$ , one has  $\text{auxd}'(l) = \text{auxd}_\Phi(l) \cap Q$ ;
- $\text{bc}' : (\mathcal{C}^{\text{prom}}(\Phi) \cap L_Q) \rightarrow \mathcal{P}(\bigcup \mathcal{Cuts}(\Phi) \cup \mathcal{D})$  is a function such that, for every  $l \in \mathcal{C}^{\text{prom}}(\Phi) \cap L_Q$ , one has  $\text{bc}'(l) = \text{bc}_\Phi(l) \cap Q$ .

With reference to notation used in definition 22,  $Q'$  is  $\mathbb{C}(\Phi)$ -complete and such that, for every  $l \in \mathcal{C}(\Phi)$  and  $p \in Q$ , if  $p \in \mathcal{P}_l^{\text{aux}}(\Phi)$  then  $\mathcal{P}_l(\Phi) \subseteq Q$ , therefore  $\text{module}_{\mathbb{C}(\Phi)}(Q')$  is well-defined.

<sup>6</sup>According to definition 19, this is equivalent to  $\mathcal{W}' = \{w \in \mathcal{W}(\Phi) \mid \exists p \in w : p \notin Q\}$ .

**Remark 23.** If  $\Phi \in \mathbf{Modules}$  and  $Q \subseteq \mathcal{P}(\Phi)$  is  $\Phi$ -complete, then  $\text{module}_\Phi(Q) \in \mathbf{Modules}$ . Furthermore, if  $\Phi \in \mathbf{PseudoPPPS}$  (resp.  $\Phi \in \mathbf{PPPS}$ ) then  $\text{nohang}(\text{module}_\Phi(Q)) \in \mathbf{PseudoPPPS}$  (resp.  $\text{nohang}(\text{module}_\Phi(Q)) \in \mathbf{PPPS}$ ).

**Definition 24** (Disjoint union of modules). *Let  $\Phi, \Phi' \in \mathbf{Modules}$ :  $\Phi$  and  $\Phi'$  are disjoint if  $\mathbb{C}(\Phi)$  and  $\mathbb{C}(\Phi')$  are disjoint,  $\mathcal{I}(\Phi) \cap \mathcal{I}(\Phi') = \emptyset$  and  $\mathcal{D}(\Phi) \cap \mathcal{D}(\Phi') = \emptyset$ .<sup>7</sup>*

Let  $n \in \mathbb{N}$  and let  $\Phi_1, \dots, \Phi_n \in \mathbf{Modules}$  be pairwise disjoint. The disjoint union of  $\Phi_1, \dots, \Phi_n$  is

$$\biguplus_{i=1}^n \Phi_i = \left( \biguplus_{i=1}^n \mathbb{C}(\Phi_i), \bigcup_{i=1}^n \mathcal{I}(\Phi_i), \bigcup_{i=1}^n \mathcal{D}(\Phi_i), \bigcup_{i=1}^n \mathcal{W}(\Phi_i), \bigcup_{i=1}^n \text{auxd}_{\Phi_i}, \bigcup_{i=1}^n \text{bc}_{\Phi_i} \right).$$

If  $n = 2$ , the disjoint union of  $\Phi_1$  and  $\Phi_2$  is denoted by  $\Phi_1 \uplus \Phi_2$ .

**Remark 25.** For every  $n \in \mathbb{N}$ , if  $\Phi_1, \dots, \Phi_n \in \mathbf{Modules}$  (resp.  $\Phi_1, \dots, \Phi_n \in \mathbf{PseudoPPPS}$ ;  $\Phi_1, \dots, \Phi_n \in \mathbf{PPPS}$ ) are pairwise disjoint and  $\Phi = \biguplus_{i=1}^n \Phi_i$ , then  $\Phi \in \mathbf{Modules}$  (resp.  $\Phi \in \mathbf{PseudoPPPS}$ ;  $\Phi \in \mathbf{PPPS}$ ).

We introduce the notion of “identity” (or better said isomorphism) between two modules. The idea is that two modules are isomorphic iff they are identical up to the names of their cells and ports (in particular, they have the same graphical representation).

**Definition 26** (Isomorphism on modules). *Let  $\Phi, \Phi' \in \mathbf{PPPS}$ .*

An isomorphism from  $\Phi$  to  $\Phi'$  is a pair  $\varphi = (\varphi_{\mathcal{C}}, \varphi_{\mathcal{P}})$  such that:

- $\varphi_{\mathcal{P}} : \mathcal{P}(\Phi) \rightarrow \mathcal{P}(\Phi')$  is a bijection where  $\text{im}(\varphi_{\mathcal{P}} \upharpoonright_{\mathcal{I}(\Phi)}) = \mathcal{I}(\Phi')$  and  $\text{im}(\varphi_{\mathcal{P}} \upharpoonright_{\mathcal{D}(\Phi)}) = \mathcal{D}(\Phi')$ ;
- $(\varphi_{\mathcal{C}}, \varphi_{\mathcal{P}} \upharpoonright_{\mathcal{P}(\mathbb{C}(\Phi))}) : \mathbb{C}(\Phi) \simeq \mathbb{C}(\Phi')$ ;
- for every  $\{p, q\} \in \mathcal{P}_2(\mathcal{P}(\Phi))$ , we have  $\{p, q\} \in \mathcal{W}(\Phi)$  iff  $\{\varphi_{\mathcal{P}}(p), \varphi_{\mathcal{P}}(q)\} \in \mathcal{W}(\Phi')$ ;
- $\text{im}(\varphi_{\mathcal{C}} \upharpoonright_{\mathcal{C}^{\text{prom}}(\Phi)}) = \mathcal{C}^{\text{prom}}(\Phi')$ ;
- the following diagrams commute:

$$\begin{array}{ccc} \mathcal{C}^{\text{prom}}(\Phi) & \xrightarrow{\text{auxd}_\Phi} & \mathcal{P}(\text{Auxdoors}(\Phi)) \\ \varphi_{\mathcal{C}} \downarrow & & \downarrow \mathcal{P}(\varphi_{\mathcal{P}}) \\ \mathcal{C}^{\text{prom}}(\Phi') & \xrightarrow{\text{auxd}_{\Phi'}} & \mathcal{P}(\text{Auxdoors}(\Phi')) \end{array} \quad \begin{array}{ccc} \mathcal{C}^{\text{prom}}(\Phi) & \xrightarrow{\text{bc}_\Phi} & \text{im}(\text{bc}_\Phi) \\ \varphi_{\mathcal{C}} \downarrow & & \downarrow \mathcal{P}(\varphi_{\mathcal{P}}) \\ \mathcal{C}^{\text{prom}}(\Phi') & \xrightarrow{\text{bc}_{\Phi'}} & \text{im}(\text{bc}_{\Phi'}) \end{array}$$

We write then  $\varphi : \Phi \simeq \Phi'$ .

If there exists an isomorphism from  $\Phi$  to  $\Phi'$ , then we say that  $\Phi$  and  $\Phi'$  are isomorphic and we write  $\Phi \simeq \Phi'$ .

<sup>7</sup>This implies that  $\mathcal{W}(\Phi) \cap \mathcal{W}(\Phi') = \emptyset$  and  $\mathcal{C}^{\text{prom}}(\Phi) \cap \mathcal{C}^{\text{prom}}(\Phi') = \emptyset$ .

**Remark 27.** Let  $\Phi, \Phi' \in \mathbf{Modules}$ . If  $\varphi$  is an isomorphism from  $\Phi$  to  $\Phi'$  then:

1.  $\text{im}(\varphi_{\mathcal{P}} \upharpoonright_{\mathcal{P}^{\text{free}}(\Phi)} = \mathcal{P}^{\text{free}}(\Phi'))$ ;
2. if  $\Phi \in \mathbf{PseudoPPPS}$  (resp.  $\Phi \in \mathbf{PPPS}$ ) then  $\Phi' \in \mathbf{PseudoPPPS}$  (resp.  $\Phi' \in \mathbf{PPPS}$ ).

### 1.3 Paths

In this section we introduce some usual notions of graph theory, merely adapted to our syntax. The only originality in our definition is that we consider the arrows (i.e. the pairs of ports connected by functions  $\text{auxd}_{\Phi}$  or  $\text{bc}_{\Phi}$ ) as edges in our undirected (hyper-)graphs.

**Definition 28** (Path, connection, acyclicity). *Let  $\Phi \in \mathbf{Modules}$ .*

*A path in  $\Phi$  is a sequence  $(p_i)_{i \in I}$  of ports of  $\Phi$  where  $I$  is an initial segment of  $\mathbb{N}$  and such that, for every  $i, i + 1 \in I$ :*

- $p_i \neq p_{i+1}$ ;
- one of the following conditions holds
  - either  $\{p_i, p_{i+1}\} \in \mathcal{W}(\Phi) \cup \text{Arrows}(\Phi)$ ,
  - or  $p_i$  and  $p_{i+1}$  are ports of a same cell of  $\Phi$ ,
- if  $i + 2 \in I$  and  $p_i = p_{i+2}$  then  $p_i$  and  $p_{i+1}$  are ports of a same cell of  $\Phi$  and  $\{p_i, p_{i+1}\} \in \mathcal{W}(\Phi)$ .<sup>8</sup>

*Let  $\varphi = (p_i)_{i \in I}$  be a path in  $\Phi$ . For every  $i \in I$ ,  $\varphi$  crosses  $p_i$ , moreover if  $i \neq 0$  and  $i + 1 \in I$  then  $\varphi$  crosses internally  $p_i$ . For every  $c \in \mathcal{W}(\Phi) \cup \text{Arrows}(\Phi)$  (resp.  $c \in \mathcal{C}(\Phi)$ ),  $\varphi$  crosses  $c$  if there exist  $i, i + 1 \in I$  such that  $\{p_i, p_{i+1}\} = c$  (resp.  $p_i, p_{i+1} \in \mathcal{P}_c(\Phi)$ ). If  $I = \emptyset$  then  $\varphi$  is the empty path. If  $I \neq \emptyset$  then  $p_0$  is the start port of  $\varphi$  (or  $\varphi$  starts from  $p_0$ ). If  $I = \{0, \dots, n\}$  for some  $n \in \mathbb{N}$ , then the path  $\varphi$  is said finite and from  $p_0$  to  $p_n$  (or connecting  $p_0$  and  $p_n$ ), furthermore  $n$  is the length of  $\varphi$  (denoted by  $\text{length}(\varphi)$ ) and  $p_n$  is the end port of  $\varphi$  (or  $\varphi$  ends in  $p_n$ ). If  $I = \mathbb{N}$ , then the path  $\varphi$  is said infinite. The terminal ports of  $\varphi$  are the start and end (if any) ports of  $\varphi$ .*

*A cycle in  $\Phi$  is a finite path  $(p_i)_{0 \leq i \leq n}$  in  $\Phi$  such that  $p_0 = p_n$  and  $n \neq 0$ .  $\Phi$  is acyclic if there is no deadlock nor cycle in  $\Phi$ .*

*Let  $p, q$  be ports of  $\rho$ :  $p$  and  $q$  are connected if there exists a path in  $\Phi$  from  $p$  to  $q$ .*

*$\Phi$  is connected if all ports  $p, q$  of  $\Phi$  are connected.*

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<sup>8</sup>This condition is imposed to facilitate the definition of cycle.

**Remark 29.** Let  $\Phi \in \mathbf{Modules}$ . If  $p \in \mathcal{I}(\Phi) \cup \mathcal{D}(\Phi)$  or  $p \in \mathcal{P}_l(\Phi)$  for some 0-ary cell  $l$  (and so  $p = \mathcal{P}_\Phi^{\text{pri}}(l)$ ) and if  $\varphi$  is a path in  $\Phi$  crossing  $p$ , then  $p$  is a terminal port of  $\varphi$ .

Notice that a finite path on a module is not empty and it might have length 0 (i.e. it consists of only one port), whereas an empty path has no length. A path on a module can cross arrows.

**Notation.** If a path in  $\Phi \in \mathbf{Modules}$  is finite, it is often denoted by a finite sequence  $(p_i)_{i \in I}$  of ports of  $\Phi$  where  $I$  is not an initial segment of  $\mathbb{N}$  but only a finite set.

The following notions will be used to compute the box associated with a promotion cell (definitions 38 and 40) in a ppps.

**Definition 30** (Ascending path, path above a promotion cell, box-crossing path). *Let  $\Phi \in \mathbf{PPPS}$ .*

*A path  $(p_i)_{i \in I}$  in  $\Phi$  (where  $I$  is an initial segment of  $\mathbb{N}$ ) is ascending if, for every  $i, i + 1 \in I$ , one of the following conditions holds:*

- *if  $p_i = \mathcal{P}_\Phi^{\text{pri}}(l)$  for some  $l \in \mathcal{C}(\Phi)$  then  $p_{i+1} \in \mathcal{P}_l^{\text{aux}}(\Phi)$ ;*
- *if  $p_i \in \mathcal{P}^{\text{aux}}(\Phi)$  then:*
  - *either  $p_{i+1} \in \mathcal{P}^{\text{pri}}(\Phi)$  with  $\{p_i, p_{i+1}\} \in \mathcal{W}(\Phi)$ ,*
  - *or  $p_{i+1} \in \text{bc}_\Phi(l) \cup \text{auxd}_\Phi(l)$  for some  $l \in \mathcal{C}^{\text{prom}}(\Phi)$  such that  $p_i \in \mathcal{P}_l^{\text{aux}}(\Phi)$ .*

*We define a binary relation  $\preceq_\Phi$  on  $\mathcal{P}(\Phi)$  by:  $p \preceq_\Phi q$  if there exists an ascending path in  $\Phi$  from  $p$  to  $q$ . For every  $n \in \mathbb{N}$ , we write  $p \preceq_\Phi^n q$  if there exists an ascending path of length  $n$  from  $p$  to  $q$ .*

*Let  $l \in \mathcal{C}^{\text{prom}}(\Phi)$ , let  $p_l$  be the (unique) auxiliary port of  $l$ . A path above<sup>9</sup>  $l$  in  $\Phi$  is an ascending path in  $\Phi$  starting from  $p_l$ . We set  $\text{cdabove}_\Phi(l) = \{q \in \bigcup \mathcal{Cuts}(\Phi) \cup \mathcal{D}(\Phi) \mid \exists \text{ path above } l \text{ ending in } q\}$ .*

*For every  $l, l' \in \mathcal{C}(\Phi)$ , we say that  $l'$  is  $\preceq$ -above<sup>9</sup>  $l$  if there exists an ascending path from  $\mathcal{P}_\Phi^{\text{pri}}(l)$  to  $\mathcal{P}_\Phi^{\text{pri}}(l')$ .*

*A path  $(p_i)_{i \in I}$  in  $\Phi$  (where  $I$  is an initial segment of  $\mathbb{N}$ ) is box-crossing if it is ascending and, for every  $i + 1 \in I$ , if  $p_{i+1} \in \text{Auxdoors}(\Phi)$  then  $p_i = \mathcal{P}_\Phi^{\text{pri}}(l)$  for the  $l \in \mathcal{C}^?(\Phi)$  such that  $p_{i+1} \in \mathcal{P}_l^{\text{aux}}(\Phi)$ .*

*We define a binary relation  $\preceq_\Phi$  on  $\mathcal{P}(\Phi)$  by:  $p \preceq_\Phi q$  if there exists a box-crossing path from  $p$  to  $q$ . For every  $n \in \mathbb{N}$ , we write  $p \preceq_\Phi^n q$  if there exists a box-crossing path of length  $n$  from  $p$  to  $q$ .*

<sup>9</sup>The use of the term ‘‘above’’ will be justified in the case of proof-structures by proposition 47.1

Roughly speaking, a box-crossing path in a ppps is an ascending path that cannot cross an arrow from a promotion cell to one of the auxiliary doors associated with it (but it can cross an arrow from a promotion cell to a cut port or a deadlock associated with it).

**Remark 31.** Let  $\Phi \in \mathbf{PPPS}$ .

1. Clearly,  $\preceq_{\Phi}^1 \subseteq \preceq_{\Phi}^1$  and  $\preceq_{\Phi} \subseteq \preceq_{\Phi}$ , furthermore  $\preceq_{\Phi} = \bigcup_{n \in \mathbb{N}} \preceq_{\Phi}^n$  and  $\preceq_{\Phi} = \bigcup_{n \in \mathbb{N}} \preceq_{\Phi}^n$ .
2. An ascending path in  $\Phi$  starting from or ending in  $p \in \mathcal{I}(\Phi) \cup \mathcal{D}_0(\Phi)$  is necessarily of length 0. More generally, if  $\varphi = (p_i)_{i \in I}$  is an ascending path in  $\Phi$  and  $p_j \in \bigcup \mathcal{Ax}(\Phi) \cup \mathcal{D}(\Phi)$  for some  $j \in I$  then  $I = \{0, \dots, j\}$  i.e.  $p_j$  is the end port of  $\varphi$ .
3.  $\preceq_{\Phi}, \preceq_{\Phi} \subseteq \mathcal{P}(\Phi)^2$  are pre-order relations, but in general they are not order relations because they are not antisymmetric. Some examples of a ppps  $\Phi$  such that  $\preceq_{\Phi}$  or  $\preceq_{\Phi}$  is not antisymmetric are given in remarks 33.2-3.
4. An ascending path  $(p_i)_{i \in I}$  in  $\Phi$  is necessarily such that if  $i + 1 \in I$  and  $p_{i+1} \in \mathbf{bc}_{\Phi}(l)$  for some  $l \in \mathcal{C}^{\text{prom}}(\Phi)$  then  $p_i$  is the unique auxiliary port of  $l$ .

## 1.4 Pre-proof-structures

Given a ppps, one might expect that there is an intuitive notion of “above/below” for its ports as done in the following definition 32, and that an axiom port is “above” a unique conclusion or cut port. But for general ppps this is wrong, because a ppps might have a “vicious cycle”, i.e. two ports which are “above” each other.

**Definition 32.** For every  $\Phi \in \mathbf{PPPS}$ , we define a the binary relation  $<_{\Phi}^1$  on  $\mathcal{P}(\Phi)$  as follows:  $p <_{\Phi}^1 p'$  if one of the following conditions holds:

- there exists a cell  $l$  of  $\Phi$  such that  $p$  is the principal port of  $l$  and  $p'$  is an auxiliary port of  $l$ ,
- $p'$  is the principal port of some cell  $l'$  of  $\Phi$ ,  $p$  is an auxiliary port of some cell  $l$  of  $\Phi$  and  $\{p, p'\}$  is a wire of  $\Phi$ .

The binary relation  $\leq_{\Phi}$  (resp.  $<_{\Phi}$ ) on  $\mathcal{P}(\Phi)$  is the reflexive-transitive (resp. transitive) closure of  $<_{\Phi}^1$ . For every  $n \in \mathbb{N}$  and  $p, p' \in \mathcal{P}(\Phi)$ , we write that  $p \leq_{\Phi}^n p'$  if there exists a finite sequence  $(p_i)_{0 \leq i \leq n}$  of ports of  $\Phi$  such that  $p_0 = p$ ,  $p_n = p'$  and  $p_i <_{\Phi}^1 p_{i+1}$  for every  $0 \leq i \leq n - 1$ .

Our definition of  $\leq_{\Phi}$  for a ppps  $\Phi$  is identical to that one in [dCT12], where we consider cut ports as minimal elements (i.e. as conclusions of  $\Phi$ ).

**Remark 33.** Let  $\Phi \in \mathbf{PPPS}$ .

1. If  $p, q \in \mathcal{P}(\Phi)$  are such that  $p \leq_{\Phi} q$  then there exists a (box-crossing) path in  $\Phi$  from  $p$  to  $q$  crossing no cuts nor axioms nor arrows. In particular, if  $q \in \mathcal{I}(\Phi) \cup \mathcal{D}(\Phi)$  then there is no  $p \in \mathcal{P}(\Phi)$  such that  $p <_{\Phi}^1 q$ .
2.  $\leq_{\Phi} \subseteq \mathcal{P}(\Phi)^2$  is a pre-order relation by definition, but in general,  $\leq_{\Phi}$  is not an order relation because it is not antisymmetric. For instance, take  $\Phi \in \mathbf{PPPS}$  consisting of a cell  $l$  such that  $\mathbf{a}_{\Phi}(l) > 0$  and a wire  $\{p, q\}$  where  $p$  is the principal port of  $l$  and  $q$  is an auxiliary port of  $l$ :  $p \leq_{\Phi} q$  (by the first condition) and  $q \leq_{\Phi} p$  (by the second condition), but  $p \neq q$ . This is an example of “vicious cycle”. A more general example of  $\Phi \in \mathbf{PPPS}$  with a “vicious cycle” is a finite sequence of cells  $l_0, \dots, l_n$  and a finite sequence of wires  $w_0, \dots, w_n$  with  $n \in \mathbb{N}$  such that  $\mathbf{a}_{\Phi}(l_i) > 0$  (where  $p_i$  and  $q_i$  are respectively the principal and an auxiliary port of  $l_i$ ) for every  $0 \leq i \leq n$ ,  $w_i = \{p_i, q_{i+1}\}$  for every  $0 \leq i \leq n-1$  and  $w_n = \{p_n, q_0\}$ .

The non-antisymmetry of  $\leq_{\Phi}$  means that if  $p, q \in \mathcal{P}(\Phi)$  are such that  $p \leq_{\Phi}^n q$  with  $n > 1$ , not necessarily  $p \neq q$ .

3. It is immediate to verify that  $<_{\Phi}^1 \subseteq \preceq_{\Phi}^1 \subseteq \preccurlyeq_{\Phi}^1$  and so  $\leq_{\Phi} \subseteq \preceq_{\Phi} \subseteq \preccurlyeq_{\Phi}$ . Therefore:
  - if  $\preccurlyeq_{\Phi}$  is antisymmetric then  $\preceq_{\Phi}$  is so;
  - if  $\preceq_{\Phi}$  is antisymmetric then  $\leq_{\Phi}$  is so.

The converses fail to hold: take for instance a ppps  $\Phi$  consisting of a 1-cell whose principal port is connected by a wire to the auxiliary port of a promotion cell  $l$  whose principal port is connected to the auxiliary port  $p$  of an unary ?-cell and such that  $\mathbf{aux}_{\Phi}(l) = \{p\}$ , then  $\leq_{\Phi}$  and  $\preceq_{\Phi}$  are antisymmetric but  $\preccurlyeq_{\Phi}$  is not.

Another example is a ppps  $\Phi$  consisting of two 1-cells whose principal ports are connected by two wires to the auxiliary ports of respectively an unary ?-cell  $l'$  and a promotion cell  $l$  such that the principal ports of  $l$  and  $l'$  are connected by a cut and  $\mathbf{bc}_{\Phi}(l) = \{\mathbf{P}_{\Phi}^{\text{pri}}(l), \mathbf{P}_{\Phi}^{\text{pri}}(l')\}$ : thus  $\leq_{\Phi}$  is antisymmetric but  $\preceq_{\Phi}$  and  $\preccurlyeq_{\Phi}$  are not.

The following lemmas 34 and 37 about the relation  $\leq_{\Phi}$  are reformulations of lemmas 10 and 14 in [dCT12] to the case of ppps with cuts.

**Lemma 34.** *Let  $\Phi \in \mathbf{PPPS}$ , let  $p, q_1, q_2 \in \mathcal{P}(\Phi)$ , let  $c, c' \in \mathcal{P}^{\text{free}}(\Phi) \cup \bigcup \mathcal{Cuts}(\Phi) \cup \mathcal{D}(\Phi)$  and let  $a \in \bigcup \mathcal{Ax}(\Phi) \cup \bigcup_{l \in \mathcal{C}(\Phi)} \{\mathbf{P}_{\Phi}^{\text{pri}}(l) \mid \mathbf{a}_{\Phi}(l) = 0\} \cup \mathcal{D}(\Phi)$ .*

1. *If  $q_1 <_{\Phi}^1 p$  and  $q_2 <_{\Phi}^1 p$  then  $q_1 = q_2$ .*

2. If  $q_1 \leq_{\Phi} p$  and  $q_2 \leq_{\Phi} p$  then  $q_1 \leq q_2$  or  $q_2 \leq q_1$ .
3. If  $p \leq_{\Phi} c$  (resp.  $a \leq_{\Phi} p$ ) then  $p = c$  (resp.  $a = p$ ).
4. If  $c \leq_{\Phi} p$  and  $c' \leq_{\Phi} p$  then  $c = c'$ .
5. If there is no  $q \in \mathcal{P}(\Phi)$  such that  $q <_{\Phi}^1 p$  (resp.  $p <_{\Phi}^1 q$ ), then  $p \in \mathcal{P}^{\text{free}}(\Phi) \cup \bigcup \mathcal{Cuts}(\Phi) \cup \mathcal{D}(\Phi)$  (resp.  $p \in \bigcup \mathcal{Ax}(\Phi) \cup \bigcup_{l \in \mathcal{C}(\Phi)} \{\mathcal{P}_{\Phi}^{\text{pri}}(l) \mid \mathbf{a}_{\Phi}(l) = 0\} \cup \mathcal{D}(\Phi)$ ).

PROOF.

1. Since every cell has exactly one principal port and because three ports cannot be connected by two wires.
2. Proof by induction on  $n \in \mathbb{N}$  where  $n$  is such that  $q_1 \leq_{\Phi}^n p$ . If  $n = 0$  then  $q_1 = p$  and so  $q_2 \leq q_1$ . If  $n > 0$  then there exists  $p_1 \in \mathcal{P}(\Phi)$  such that  $q_1 \leq_{\Phi}^{n-1} p_1 <_{\Phi}^1 p$ : if  $q_2 = p$  then  $q_1 \leq q_2$ ; otherwise there exists  $p_2 \in \mathcal{P}(\Phi)$  such that  $q_2 \leq_{\Phi}^n p_2 <_{\Phi}^1 p$ , so  $p_1 = p_2$  by lemma 34.1, therefore  $q_1 \leq_{\Phi} q_2$  or  $q_2 \leq_{\Phi} q_1$  by induction hypothesis applied to  $p_1$ .
3. If  $c \in \mathcal{P}^{\text{free}}(\Phi) \cup \bigcup \mathcal{Cuts}(\Phi) \cup \mathcal{D}(\Phi)$  (resp.  $a \in \bigcup \mathcal{Ax}(\Phi) \cup \bigcup_{l \in \mathcal{C}(\Phi)} \{\mathcal{P}_{\Phi}^{\text{pri}}(l) \mid \mathbf{a}_{\Phi}(l) = 0\} \cup \mathcal{D}(\Phi)$ ), then  $p \not<_{\Phi}^1 c$  (resp.  $a \not<_{\Phi}^1 p$ ) for every  $p \in \mathcal{P}(\Phi)$ .
4. By lemma 34.2,  $c \leq c'$  or  $c' \leq c$ ; in any case,  $c = c'$  by lemma 34.3.
5. If  $p \notin \mathcal{P}^{\text{free}}(\Phi) \cup \bigcup \mathcal{Cuts}(\Phi) \cup \mathcal{D}(\Phi)$  (resp.  $p \notin \bigcup \mathcal{Ax}(\Phi) \cup \bigcup_{l \in \mathcal{C}(\Phi)} \{\mathcal{P}_{\Phi}^{\text{pri}}(l) \mid \mathbf{a}_{\Phi}(l) = 0\} \cup \mathcal{D}(\Phi)$ ), then there are only two cases:
  - either  $p \in \mathcal{P}^{\text{pri}}(\Phi)$  (resp.  $p \in \mathcal{P}^{\text{aux}}(\Phi)$ ) and there exists  $q \in \mathcal{P}^{\text{aux}}(\Phi)$  (resp.  $q \in \mathcal{P}^{\text{pri}}(\Phi)$ ) such that  $\{p, q\} \in \mathcal{W}(\Phi)$ , so  $q <_{\Phi}^1 p$  (resp.  $p <_{\Phi}^1 q$ );
  - or  $p \in \mathcal{P}_l^{\text{aux}}(\Phi)$  (resp.  $p = \mathcal{P}_{\Phi}^{\text{pri}}(l)$ ) for some cell  $l$  of  $\Phi$  such that  $\mathcal{P}_l^{\text{aux}}(\Phi) \neq \emptyset$ , so  $\mathcal{P}_{\Phi}^{\text{pri}}(l) <_{\Phi}^1 p$  (resp. there exists  $q \in \mathcal{P}_l^{\text{aux}}(\Phi)$  such that  $p <_{\Phi}^1 q$ ).

□

Lemma 34.3 means that conclusions, cuts ports and deadlocks (resp. axiom ports, principal ports of 0-ary cells and deadlocks) of a ppps  $\Phi$  are the minimal (resp. maximal) elements of the pre-order relation  $\leq_{\Phi}$ . Lemma 34.4 implies that in a ppps  $\Phi$ , an axiom port cannot be “above” (in the sense of definition 32) two different conclusions or cut ports of  $\Phi$ .

**Definition 35** (Pre-proof-structure). *A pre-proof-structure (or pps for short) is a  $\Phi \in \mathbf{PPPS}$  such that  $\leq_{\Phi}$  is antisymmetric.*

*We denote by  $\mathbf{PPS}$  the set of pre-proof-structures.*

*We denote by  $\mathbf{PPS}_{\text{DiLL}_0}$  the set of  $\Phi \in \mathbf{PPS}$  such that  $\mathcal{C}^{\text{prom}}(\Phi) = \emptyset$ , whose elements are the DiLL<sub>0</sub>-proof-structures<sup>10</sup> (or DiLL<sub>0</sub>-pps for short).*

<sup>10</sup>We have deliberately forgotten a “pre-”. The reason will explained in remark 45.



We remind that in a pps  $\Phi$ , antisymmetry of the relation  $\leq_\Phi$  entails that  $\leq_\Phi$  is an order: this prevents from creating in  $\Phi$  the “vicious cycles” seen in remark 33.2. Therefore, we can see  $\leq_\Phi$  as a definition of a relation “above/below” for the ports of a pps  $\Phi$ . We will see that  $\preceq_\Phi$  and  $\preccurlyeq_\Phi$  extend this relation in the case of a proof-structure  $\Phi$  (see proposition 47.1).

**Remark 36.** Let  $\Phi \in \mathbf{PPS}$  and  $\Phi' \in \mathbf{PPPS}$ . If  $\Phi \simeq \Phi'$  then  $\Phi' \in \mathbf{PPS}$ .

**Lemma 37.** Let  $\Phi \in \mathbf{PPS}$  and let  $p \in \mathcal{P}(\Phi)$ :

1. there exists exactly one  $c \in \mathcal{P}^{\text{free}}(\Phi) \cup \bigcup \mathcal{Cuts}(\Phi) \cup \mathcal{D}(\Phi)$  such that  $c \leq_\Phi p$ .
2. there exists at least one  $a \in \bigcup \mathcal{Ax}(\Phi) \cup \bigcup_{l \in \mathcal{C}(\Phi)} \{\mathbf{P}_\Phi^{\text{pri}}(l) \mid \mathbf{a}_\Phi(l) = 0\} \cup \mathcal{D}(\Phi)$  such that  $p \leq_\Phi a$ .

PROOF.

1. For the unicity, apply lemma 34.4. For the existence, we build an initial segment  $I$  of  $\mathbb{N}$  and a non-empty “downward” path  $(p_i)_{i \in I}$  in  $\Phi$  as follows:

- $0 \in I$  and  $p_0 = p$ ;
- if  $i \in I$  then:
  - if  $p_i \in \mathcal{P}^{\text{free}}(\Phi) \cup \bigcup \mathcal{Cuts}(\Phi) \cup \mathcal{D}(\Phi)$  then  $I = \{0, \dots, i\}$ ,
  - otherwise, by lemma 34.5, there exists  $q \in \mathcal{P}(\Phi)$  such that  $q <_\Phi^1 p$ , and then  $i + 1 \in I$  and  $p_{i+1} = q$ .

By construction,  $p_{i+1} <_\Phi^1 p_i$  for every  $i, i + 1 \in I$ . By antisymmetry of  $\leq_\Phi$ ,  $I$  is finite (otherwise there would be a “vicious cycle” as  $\mathcal{P}(\Phi)$  is a finite set, i.e. there would exist  $i, j \in I$  such that  $i < j$  and  $p_i = p_j$ , so  $p_i <_\Phi^1 p_{j-1}$  and  $p_{j-1} \leq_\Phi p_i$ , that is impossible by antisymmetry of  $\leq_\Phi$ ), hence there exists  $n \in \mathbb{N}$  such that  $I = \{0, \dots, n\}$  and  $p_n \in \mathcal{P}^{\text{free}}(\Phi) \cup \bigcup \mathcal{Cuts}(\Phi) \cup \mathcal{D}(\Phi)$ , with  $p_n \leq_\Phi p$ .

2. We build an initial segment  $I$  of  $\mathbb{N}$  and a non-empty “upward” path  $(p_i)_{i \in I}$  in  $\Phi$  as follows:

- $0 \in I$  and  $p_0 = p$ ;
- if  $i \in I$  then:
  - if  $p_i \in \bigcup \mathcal{Ax}(\Phi) \cup \bigcup_{l \in \mathcal{C}(\Phi)} \{\mathbf{P}_\Phi^{\text{pri}}(l) \mid \mathbf{a}_\Phi(l) = 0\} \cup \mathcal{D}(\Phi)$  then  $I = \{0, \dots, i\}$ ,
  - otherwise, by lemma 34.5, there exists  $q \in \mathcal{P}(\Phi)$  such that  $p <_\Phi^1 q$ , and then  $i + 1 \in I$  and  $p_{i+1} = q$ .

By construction,  $p_i <_{\Phi}^1 p_{i+1}$  for every  $i, i+1 \in I$ . By antisymmetry of  $\leq_{\Phi}$ ,  $I$  is finite (otherwise there would be a “vicious cycle” as  $\mathcal{P}(\Phi)$  is a finite set, i.e. there would exist  $i, j \in I$  such that  $i < j$  and  $p_i = p_j$ , so  $p_{j-1} <_{\Phi}^1 p_i$  and  $p_i \leq_{\Phi} p_{j-1}$ , that is impossible by antisymmetry of  $\leq_{\Phi}$ ), hence there exists  $n \in \mathbb{N}$  such that  $I = \{0, \dots, n\}$  and  $p_n \in \bigcup \mathcal{A}x(\Phi) \cup \bigcup_{l \in \mathcal{C}(\Phi)} \{P_{\Phi}^{\text{pri}}(l) \mid \mathbf{a}_{\Phi}(l) = 0\} \cup \mathcal{D}(\Phi)$ , with  $p \leq_{\Phi} p_n$ .  $\square$

Lemma 37.1 entails that in a pps  $\Phi$ , any axiom port is above a unique conclusion or cut port of  $\Phi$ . Lemmas 37.1-2 means that if  $\Phi$  is a pps then the order relation  $\leq_{\Phi}$  defines a natural “top-down” orientation on ports from axiom ports and principal ports of 0-ary cells to conclusions and cut ports of  $\Phi$ , where by lemma 34.3 deadlocks of  $\Phi$  have no ports above or below them (in the sense of definition 32).

## 1.5 Boxes and (non-inductive) proof-structures

Similarly to [dCT12], the main difference of our syntax from the usual syntaxes of linear logic (or differential linear logic with boxes) proof-nets (see for example [Gir87, Lau03, Pag09, dCPT11, Tra11]) is the absence of an explicit (inductive) constructor for boxes: this leads to define a box as a sort of sub-graph satisfying some conditions. This more “geometrical” approach was followed for example in [DR95, Tor03, MP07, dCT12]. In our syntax we have to reconstruct the boxes of a pps  $\Phi$  by using some “geometrical” informations coming from  $\Phi$ , in particular the arrow functions  $\mathbf{auxd}_{\Phi}$  and  $\mathbf{bc}_{\Phi}$  play a crucial role. Each promotion cell in a pre-proof-structure corresponds to the so-called “principal ports of a box” in the usual syntaxes of linear logic proof-nets. More delicate is the issue of marking out the other boundaries of a box, which are called “auxiliary ports of a box” in the usual syntax (corresponding to auxiliary doors in our syntax), and the content of a box: in order to do that, some conditions are to be fulfilled.

**Definition 38** (Box). *Let  $\Phi \in \mathbf{PPS}$ .*

*Let  $l \in \mathcal{C}^{\text{prom}}(\Phi)$  and let  $p_l$  be the unique auxiliary port of  $l$ .<sup>11</sup> We say that “the box of  $l$  is defined in  $\Phi$ ” or “ $l$  has a box in  $\Phi$ ” when, for every  $q, q' \in \mathcal{P}(\Phi)$ :*

1. *if  $q \in \mathbf{auxd}_{\Phi}(l)$  then  $q \not\leq_{\Phi} p_l$  and  $p_l \not\leq_{\Phi} q$ ;*
2. *if  $q, q' \in \mathbf{auxd}_{\Phi}(l)$  with  $q \neq q'$  then  $q \not\leq_{\Phi} q'$  and  $q' \not\leq_{\Phi} q$ ;*
3. *for every  $l' \in \mathcal{C}^{\text{prom}}(\Phi) \prec\text{-above } l$ , if  $q' \in \mathbf{cutports}_{\Phi}(l')$  and  $q \in \mathbf{doors}_{\Phi}(l)$ , then  $q' \not\leq_{\Phi} q$ ;*

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<sup>11</sup>Remind that  $\mathbf{a}_{\Phi}(l) = 1$  and  $\mathbf{t}_{\Phi}(l) = !$  since  $l$  is a promotion cell.

4. if  $q$  is the end port of a path above  $l$  in  $\Phi$  and if  $q' = q$  or  $\{q, q'\} \in \mathcal{Ax}(\Phi)$ , then there exists  $r \in \text{doors}_\Phi(l) \cup \text{cdabove}_\Phi(l)$  such that  $r \leq_\Phi q'$ .<sup>12</sup>

We denote by  $\mathcal{C}^{\text{box}}(\Phi)$  the set of  $l \in \mathcal{C}^{\text{prom}}(\Phi)$  having a box in  $\Phi$ , whose elements are the box cells of  $\Phi$ .

**Remark 39.**  $\mathcal{C}^{\text{box}}(\Phi) = \emptyset$  for every  $\Phi \in \mathbf{PPS}_{\text{DiLL}_0}$ , as  $\mathcal{C}^{\text{prom}}(\Phi) = \emptyset$ . Hence  $\leq_\Phi = \preceq_\Phi = \preccurlyeq_\Phi$  for every  $\Phi \in \mathbf{PPS}_{\text{DiLL}_0}$

When a promotion cell  $l$  satisfies conditions 1, 2, 3 and 4 in definition 38, we are able to compute the box associated with  $l$ , by taking into account only the “geometrical” informations available in a pps. The following definition of how to compute a box is quite delicate and it is inspired by the analogous definition in [dCT12], with two further complications: our definition is completely non-inductive and our boxes might contain cuts. Intuitively, for every  $l \in \mathcal{C}^{\text{box}}(\Phi)$ ,  $\text{auxd}_\Phi(l)$  gives the auxiliary doors of the box  $\text{box}_\Phi(l)$  associated with  $l$  (i.e. the boundaries of this box as well as  $l$ ),  $\text{bc}_\Phi(l)$  gives the cut ports and deadlocks belonging to  $\text{box}_\Phi(l)$  and not belonging to more inner boxes, the content of  $\text{box}_\Phi(l)$  is “all that is above” the doors of  $\text{box}_\Phi(l)$  and the cut ports pointed by  $\text{bc}_\Phi(l)$ .

**Definition 40** (Computation of a box). *Let  $\Phi \in \mathbf{PPS}$ , let  $v \in \mathcal{C}^{\text{box}}(\Phi)$  and let  $r_v$  be the unique auxiliary port of  $v$ . We set:*

$$\begin{aligned} \text{inbox}_\Phi(v) &= \{p \in \mathcal{P}(\Phi) \mid \exists \text{ path above } v \text{ in } \Phi \text{ ending in } p\} \\ B'_v &= \text{inbox}_\Phi(v) \setminus \{r_v\} \\ B_v &= \begin{cases} \text{inbox}_\Phi(v) & \text{if } r_v \in \bigcup \mathcal{Ax}(\Phi) \\ B'_v & \text{otherwise.} \end{cases} \end{aligned}$$

$\text{inbox}_\Phi(v)$  is the content of the box of  $v$  in  $\Phi$ .

Let  $\mathcal{L}_0$  and  $\mathcal{P}_0$  be two sets such that there exist two bijections  $p_1 : \mathcal{L}_0 \rightarrow \text{auxd}_\Phi(v)$  and  $p_0 : \mathcal{L}_0 \rightarrow \mathcal{P}_0$ , moreover  $\mathcal{L}_0 \cap (\mathcal{P}(\mathcal{C}_\Phi)(B'_v) \setminus \mathcal{P}(\mathcal{C}_\Phi)(\text{auxd}_\Phi(v))) = \emptyset$  and  $\mathcal{P}_0 \cap B_v = \emptyset$ . We set  $\mathcal{C}_v = (\mathfrak{t}_v, \mathcal{P}_v, \mathcal{C}_v, \mathbf{P}_v^{\text{pri}}, \mathbf{P}_v^{\text{left}})$  where:

- $\mathfrak{t}_v$  is a function from  $\mathcal{C}(\mathcal{C}_v) = \mathcal{L}_0 \cup (\mathcal{P}(\mathcal{C}_\Phi)(B'_v) \setminus \mathcal{P}(\mathcal{C}_\Phi)(\text{auxd}_\Phi(v)))$  to  $\mathcal{T}$  such that  $\mathfrak{t}_v(l) = ?$  for every  $l \in \mathcal{L}_0$  and  $\mathfrak{t}_v \upharpoonright_{\mathcal{P}(\mathcal{C}_\Phi)(B'_v) \setminus \mathcal{P}(\mathcal{C}_\Phi)(\text{auxd}_\Phi(v))} = \mathfrak{t}_\Phi \upharpoonright_{\mathcal{P}(\mathcal{C}_\Phi)(B'_v) \setminus \mathcal{P}(\mathcal{C}_\Phi)(\text{auxd}_\Phi(v))}$ ;<sup>13</sup>
- $\mathcal{P}_v = \mathcal{P}_0 \cup B'_v$ ;<sup>14</sup>

<sup>12</sup>Because of the definition of  $\text{auxd}_\Phi$  and  $\text{cdabove}_\Phi$ , this condition means that if  $l$  has a box in  $\Phi$  then either  $r = p_i$  or  $\text{auxd}_\Phi$  has to point from  $l$  to  $r$  or  $\text{bc}_\Phi$  has to point from  $l'$  to  $r$ , where  $l' \in \mathcal{C}^{\text{prom}}(\Phi)$  is “above”  $l$  (in the sense of definition 30). In particular, there exists a path above  $l$  ending in  $q'$ .

<sup>13</sup>Notice that  $v = \mathcal{C}_\Phi(r_v) \notin \mathcal{C}(\mathcal{C}_v)$ .

<sup>14</sup>Note that  $r_v \notin \mathcal{P}_v$ .

- $\mathbb{C}_v : \mathcal{P}_v \rightarrow \mathcal{C}(\mathbb{C}_v)$  is such that, for every  $p \in \mathcal{P}_v$ ,

$$\mathbb{C}_v(p) = \begin{cases} \mathbb{C}_\Phi(p) & \text{if } p \in B'_v \setminus \text{auxd}_\Phi(v) \\ l & \text{if } p = p_0(l) \text{ or } p = p_1(l) \text{ for some } l \in \mathcal{L}_0 \end{cases}$$

- $\mathbb{P}_v^{\text{pri}} : \mathcal{C}(\mathbb{C}_v) \rightarrow \mathcal{P}_v$  is such that, for every  $l \in \mathcal{C}(\mathbb{C}_v)$ ,

$$\mathbb{P}_v^{\text{pri}}(l) = \begin{cases} p_0(l) & \text{if } l \in \mathcal{L}_0 \\ \mathbb{P}_\Phi^{\text{pri}}(l) & \text{otherwise} \end{cases}$$

- $\mathbb{P}_v^{\text{left}} = \mathbb{P}_\Phi^{\text{left}} \upharpoonright_{\mathcal{C}^{\otimes, \mathfrak{A}}(\Phi) \cap \mathcal{C}(\mathbb{C}_v)}$ .

The box of  $v$  in  $\Phi$  is  $\text{box}_\Phi(v) = (\mathbb{C}_v, \mathcal{I}_v, \mathcal{D}_v, \mathcal{W}_v, \text{auxd}_v, \text{bc}_v)$  where:

- 

$$\mathcal{I}_v = \begin{cases} \{r_v\} & \text{if } r_v \in B_v \\ \emptyset & \text{otherwise;} \end{cases}$$

- $\mathcal{W}_v = \{w \in \mathcal{W}(\Phi) \mid w \subseteq B_v\}$ ;
- $\mathcal{D}_v = \mathcal{D}(\Phi) \cap \text{inbox}_\Phi(v)$ ;
- $\text{auxd}_v = \text{auxd}_\Phi \upharpoonright_{\mathcal{C}^{\text{prom}}(\Phi) \cap \mathcal{C}(\mathbb{C}_v)}$ ;
- $\text{bc}_v = \text{bc}_\Phi \upharpoonright_{\mathcal{C}^{\text{prom}}(\Phi) \cap \mathcal{C}(\mathbb{C}_v)}$ .

The idea in definition 40 is that, given  $\Phi \in \mathbf{PPS}$ , we build the box  $\text{box}_\Phi(v)$  of  $v \in \mathcal{C}^{\text{box}}(\Phi)$  by starting from  $v$  by means of the paths above  $v$  (to get  $\text{inbox}_\Phi(v)$ ) and then by reconstructing  $\text{box}_\Phi(v)$  as the (in some sense) “smallest sub-pps” of  $\Phi$  containing  $\text{inbox}_\Phi(v)$ : nevertheless even if we give a precise definition of sub-pps, it is not correct to say that  $\text{box}_\Phi(v)$  is the smallest sub-pps of  $\Phi$  containing  $\text{inbox}_\Phi(v)$  because  $\text{inbox}_\Phi(v) \subseteq \mathcal{P}(\Phi)$  but in general  $\mathcal{P}_v \not\subseteq \mathcal{P}(\Phi)$ . Some syntactical complications in the definition of  $\text{box}_\Phi(v)$  are due to get that propositions 41 and 46 hold, for example the issue whether the auxiliary port of  $v$  belongs or not to  $\text{box}_\Phi(v)$ , and the add of the sets  $\mathcal{L}_0$  and  $\mathcal{P}_0$ : every  $l \in \mathcal{L}_0$  is a unary ?-cell, where  $p_0(l)$  (resp.  $p_1(l)$ ) is its principal (resp. unique auxiliary) port. Note that  $\text{auxd}_v$  is like  $\text{auxd}_\Phi$  but it forgets all the arrows associated with the promotion cells of  $\Phi$  that are not in  $\mathcal{C}(\mathbb{C}_v)$ , including  $v$ . Similarly for  $\text{bc}_v$ .

**Proposition 41.** *Let  $\Phi \in \mathbf{PPS}$  and let  $v \in \mathcal{C}^{\text{box}}(\Phi)$ . With reference to notation of definition 40.*

1.  $\mathbb{C}_v \in \mathbf{Cells}$  and  $l \in \mathcal{C}(\Phi)$  for every  $l \in \mathcal{C}(\mathbb{C}_v) \setminus \mathcal{L}_0$ .
2.  $\text{box}_\Phi(v) \in \mathbf{PPS}$  with  $\mathbb{C}(\text{box}_\Phi(v)) = \mathbb{C}_v$ .

PROOF. Intuitive?  $\square$

**Remark 42.** Let  $\Phi \in \mathbf{PPS}$ .

1. Given  $v \in \mathcal{C}^{\text{box}}(\Phi)$  whose unique auxiliary port is  $r_v$ , one has  $r_v \in \text{inbox}_{\Phi}(v)$  and  $\{r_v, p_v\} \in \mathcal{W}(\Phi)$  for some  $p_v \in \text{inbox}_{\Phi}(v) \cap \mathcal{P}(\text{box}_{\Phi}(v))$ . If  $\{r_v, p_v\} \in \mathcal{Ax}(\Phi)$  (resp.  $\{r_v, p_v\} \notin \mathcal{Ax}(\Phi)$ ) then  $\text{inbox}_{\Phi}(v) \subseteq \mathcal{P}(\text{box}_{\Phi}(v))$  (resp.  $\text{inbox}_{\Phi}(v) \setminus \{r_v\} \subseteq \mathcal{P}(\text{box}_{\Phi}(v))$ ),  $\{r_v, p_v\} \in \mathcal{W}(\text{box}_{\Phi}(v))$  (resp.  $\{r_v, p_v\} \notin \mathcal{W}(\text{box}_{\Phi}(v))$ ) and  $r_v$  (resp.  $p_v$ ) is the unique  $q \in \mathcal{P}^{\text{free}}(\text{box}_{\Phi}(v))$  such that  $\mathbb{P}_{\Phi}^{\text{pri}}(v) <_{\Phi} q$ .
2. Let  $v, v' \in \mathcal{C}^{\text{box}}(\Phi)$ : there exists a path above  $v$  in  $\Phi$  ending in the unique auxiliary port  $p_{v'}$  of  $v'$  iff  $\text{inbox}_R(v') \subseteq \text{inbox}_R(v)$ . Indeed, for the left-to-right direction, for every  $p \in \mathcal{P}(\Phi)$ , if  $p \in \text{inbox}_R(v')$  then there exists a path  $\varphi'$  above  $v'$  in  $\Phi$  ending in  $p$ ; by hypothesis and definition of ascending path, there exists a path  $\varphi$  above  $v$  in  $\Phi$  ending in  $\mathbb{P}_{\Phi}^{\text{pri}}(v')$ , hence  $\varphi \cdot \varphi'$  is a path above  $v$  in  $\Phi$  ending in  $p$ , hence  $p \in \text{inbox}_R(v)$ . Conversely,  $p_{v'} \in \text{inbox}_R(v')$   $\subseteq$   $\text{inbox}_R(v)$ , so there exists a path above  $v$  in  $\Phi$  ending in  $p_{v'}$ .

The following lemma shows an expected property of boxes in a pps: all that is below (in the sense of definition 32) the doors of a box cannot be inside this box, i.e. the doors of the box associated with a promotion cell are the boundaries of this box. The proof of this lemma uses all the conditions mentioned in definition 38, so it reveals indirectly their importance.

**Lemma 43.** *Let  $\Phi \in \mathbf{PPS}$  and  $v \in \mathcal{C}^{\text{box}}(\Phi)$ . For every  $p, q \in \mathcal{P}(\Phi)$ , if  $p <_{\Phi} q$  and  $q \in \text{doors}_{\Phi}(v)$  then  $p \notin \text{inbox}_{\Phi}(v)$ .*

PROOF. Let  $q \in \text{doors}_{\Phi}(v)$ . If  $p \in \text{inbox}_{\Phi}(v)$  then there exists a path above  $v$  ending in  $p$ . By condition 4 in definition 38, there exists  $r \in \text{doors}_{\Phi}(v) \cup \text{cdabove}_{\Phi}(v)$  such that  $r \leq_{\Phi} p$ . If  $r \in \text{doors}_{\Phi}(v)$  then  $p \not<_{\Phi} q$  by conditions 1-2 of definition 38. Otherwise  $r \in \text{cdabove}_{\Phi}(v)$  and then there exists  $v' \in \mathcal{C}^{\text{prom}}(\Phi) \preceq$ -above  $v$  such that  $r \in \text{bc}_{\Phi}(v')$ ; there are only two cases: either  $r \in \text{deadlocks}_{\Phi}(v')$  and so  $r = p$  and  $p \not<_{\Phi} q$  by lemma 34.3, or  $r \in \text{cutports}_{\Phi}(v')$  and so  $p \not<_{\Phi} q$  otherwise  $r \leq_{\Phi} p <_{\Phi} q$  that is impossible by condition 3 in definition 38.  $\square$

We can introduce now the syntactical objects for which we prove our main result: proof-structures.

**Definition 44** (Proof-structure). *A proof-structure (or ps for short) is a  $R \in \mathbf{PPS}$  such that:*

- $\mathcal{C}^{\text{box}}(R) = \mathcal{C}^{\text{prom}}(R)$ ,<sup>15</sup>

<sup>15</sup>This means that every promotion cell of  $R$  has a box in  $R$  (in the sense of definition 38).

- (nesting condition) for every  $l, l' \in \mathcal{C}^{\text{box}}(R)$  either  $\text{inbox}_R(l) \subseteq \text{inbox}_R(l')$  or  $\text{inbox}_R(l') \subseteq \text{inbox}_R(l)$  or  $\text{inbox}_R(l) \cap \text{inbox}_R(l') = \emptyset$ ;

We denote by  $\mathbf{PS}$  the set of proof-structures.

We denote by  $\mathbf{PS}_{\text{MELL}}$  the set of  $R \in \mathbf{PS}$  such that  $\mathcal{C}^!(R) = \mathcal{C}^{\text{box}}(R)$ , whose elements are the MELL-proof-structures (or MELL-ps for short).

We point out that our definition of proof-structure is completely non-inductive, as in [DR95, Tor03, MP07], differently from the usual definitions of proof-structure in the literature on linear logic and its differential version (see for example [Gir87, Lau03, Pag09, dCPT11, Tra11, dCT12]). This leads to consider the linear logic proof-structures as “really geometrical” objects, in accordance with the Girard’s original spirit. This definition allows also to define the cut-elimination directly on these “geometrical” objects. Actually the definition of cut-free proof-structure given in [dCT12] can be reformulated in a non-inductive way (this remark was our starting point).

**Remark 45.**  $\mathbf{PPS}_{\text{DiLL}_0} \subseteq \mathbf{PS}$ , since  $\mathcal{C}^{\text{box}}(\Phi) \subseteq \mathcal{C}^{\text{prom}}(\Phi) = \emptyset$  for every  $\Phi \in \mathbf{PPS}_{\text{DiLL}_0}$ .

**Proposition 46.** Let  $R \in \mathbf{PS}$  and let  $l \in \mathcal{C}^{\text{box}}(\Phi)$ .

1.  $\text{box}_R(l) \in \mathbf{PS}$ .
2. If  $\varphi : R \simeq R'$  then  $\varphi_{\mathcal{P}}(l) \in \mathcal{C}^{\text{box}}(R')$  and  $\text{box}_R(l) \simeq \text{box}_{R'}(\varphi_{\mathcal{C}}(l))$ .

WHY?

PROOF. Intuitive? □

The following proposition says that in a ps  $R$ ,  $\preceq_R$  and  $\preccurlyeq_R$  are order relations. In a certain sense, in the case of a ps  $R$ ,  $\preceq_R$  and  $\preccurlyeq_R$  are “good generalizations” of  $\leq_R$  (they extend  $\leq_{\Phi}$  as a relation “above/below” for the ports of  $R$ : see remark 33.3 but also the following lemma 49), with the further property that any promotion cell is the least element (with respect to  $\preccurlyeq_R$ ) of the box associated with it (proposition 47.2).

**Proposition 47.** Let  $R \in \mathbf{PS}$ .

1.  $\preceq_R$  and  $\preccurlyeq_R$  are order relations on  $\mathcal{P}(R)$ .
2. For every  $v \in \mathcal{C}^{\text{box}}(R)$ , the unique auxiliary port of  $v$  is the least element in  $\text{inbox}_R(v)$  with respect to  $\preccurlyeq_R$ .
3. For every  $v, v' \in \mathcal{C}^{\text{box}}(R)$ , if  $v \neq v'$  then  $\text{inbox}_R(v) \neq \text{inbox}_R(v')$ .

PROOF.

1. By remarks 31.3 and 33.3, it suffices to show that  $\preccurlyeq_R$  is antisymmetric. Let us suppose by absurd that  $\preccurlyeq_R$  is not antisymmetric, so there exist  $p, q \in \mathcal{P}(R)$  such that  $p \preccurlyeq_R q$ ,  $q \preccurlyeq_R p$  and  $p \neq q$ . As  $R \in \mathbf{PPS}$ ,  $\leq_R$  is

antisymmetric, thus it is impossible that  $p \leq_R q$  and  $q \leq_R p$ . Hence, there exist  $p_0, p_1, p_2 \in \mathcal{P}(R)$  and  $v \in \mathcal{C}^{\text{prom}}(R) = \mathcal{C}^{\text{box}}(R)$  (as  $R \in \mathbf{PS}$ ) such that  $p_0$  (resp.  $p_1$ ) is the principal (resp. unique auxiliary) port of  $v$ ,  $p_2 \in \text{auxd}_R(v) \cup \text{cutports}_R(v)$  and  $p_2 \preceq_R p_0$  ( $p_2 \notin \text{deadlocks}_R(v)$  by remark 31.2 since  $p_2 \preceq_R^n p_1$  for some  $n \in \mathbb{N}^*$ ). Therefore  $p_0 \in \text{inbox}_R(v)$ ,  $p_0 <_R p_1$  and  $p_1 \in \text{doors}_R(v)$ , that is impossible by lemma 43.

2. By definition of  $\text{inbox}_R(v)$  and  $\preceq_R$ , if  $p$  is the unique auxiliary port of  $v$  then  $p \preceq_R q$  for every  $q \in \text{inbox}_R(v)$ . We conclude thanks to proposition 47.1.
3. Let  $p_v$  (resp.  $p_{v'}$ ) the unique auxiliary port of  $v$  (resp.  $v'$ ). If there is not path above  $v$  ending in  $p_{v'}$  then  $p_{v'} \notin \text{inbox}_R(v)$  but  $p_{v'} \in \text{inbox}_R(v')$ , hence  $\text{inbox}_R(v) \neq \text{inbox}_R(v')$ . Otherwise there exists a path above  $v$  ending in  $p_{v'}$  (thus  $p_v \preceq_R p_{v'}$ ), moreover  $p_v \neq p_{v'}$  since  $v \neq v'$ ; by antisymmetry of  $\preceq_R$  (proposition 47.1), there is no path above  $v'$  ending in  $p_v$ , so  $p_v \notin \text{inbox}_R(v')$  but  $p_{v'} \in \text{inbox}_R(v')$ , therefore  $\text{inbox}_R(v) \neq \text{inbox}_R(v')$ . □

**Remark 48.** In the proof of proposition 47 the hypothesis that  $R \in \mathbf{PS}$  is used only to ensure that every promotion cell of  $R$  has a box defined in  $R$ , in particular we never used the hypothesis that  $R$  fulfills the nesting condition. Notice that in the examples showed in remark 33.3 the promotion cells have no box defined. An example of pps  $\Phi$  such that  $\mathcal{C}^{\text{box}}(\Phi) = \mathcal{C}^{\text{prom}}(\Phi)$  (and so  $\preceq_\Phi$  is an order relation) but the nesting condition is not fulfilled (and so  $\Phi \notin \mathbf{PS}$ ) is the following: take one  $\perp$ -cell whose principal port is connected by a wire to the auxiliary port  $p$  of an unary  $?$ -cell and two 1-cells whose principal ports are both connected by a wire respectively to the unique auxiliary ports  $p_v$  and  $p_{v'}$  of two  $!$ -cells  $v$  and  $v'$  which have both an arrow pointing to  $p$ ; in this case  $p \in \text{inbox}_\Phi(v) \cap \text{inbox}_\Phi(v')$  but  $\text{inbox}_\Phi(v) \not\subseteq \text{inbox}_\Phi(v')$  (because the  $p_v \in \text{inbox}_R(v) \setminus \text{inbox}_R(v')$ ) and  $\text{inbox}_R(v) \not\subseteq \text{inbox}_R(v')$  (because  $p_{v'} \in \text{inbox}_R(v') \setminus \text{inbox}_R(v)$ ).

Given a ps  $R$ , we can generalize lemmas 34 and 37 for the order relation  $\preceq_R$ , which is a restriction of  $\preceq_R$  (see remark 33.3).

**Lemma 49.** *Let  $R \in \mathbf{PS}$  and  $p, p', q \in \mathcal{P}(R)$ .*

1. *If  $p \preceq_R^1 q$  and  $p' \preceq_R^1 q$  then  $p = p'$ .*
2. *If  $p \preceq_R q$  and  $p' \preceq_R q'$  then either  $p \preceq_R p'$  or  $p' \preceq_R p$ .*
3. *For every  $c \in \mathcal{P}^{\text{free}}(R) \cup \bigcup \text{Cuts}_0(R) \cup \mathcal{D}_0(R)$ , if there exists an ascending path from  $q$  to  $c$  then  $q = c$ .*
4. *There exists at most one box-crossing path in  $R$  from  $p$  to  $q$ .*

5. There exists a unique  $c \in \mathcal{P}^{\text{free}}(R) \cup \bigcup \mathcal{Cuts}_0(R) \cup \mathcal{D}_0(R)$  such that  $c \preceq_R q$ .

PROOF.

1.  $q \notin \mathcal{I}(R) \cup \mathcal{D}_0(R)$  by remark 31.2, hence there are only three cases:
  - $q \in \mathcal{P}^{\text{pri}}(R)$ , so  $p \in \mathcal{P}^{\text{aux}}(R)$  and either  $\{p, q\} \in \mathcal{W}(R)$  or  $p$  is the unique auxiliary port of some  $l \in \mathcal{C}^{\text{box}}(R)$  such that  $q \in \text{cutports}_R(l)$ ; analogously for  $p'$ ; by definition of ppps (in particular, condition 1 about the set of wires in definition 12), necessarily  $p = p'$ ;
  - $q \in \mathcal{P}^{\text{aux}}(R)$ , so  $p = \mathbf{P}_R^{\text{pri}}(l) = p'$  for the  $l \in \mathcal{C}(R)$  such that  $q \in \mathcal{P}_l^{\text{aux}}(R)$  (because of the definition of box-crossing path);
  - $q \in \mathcal{D}(R)$  and there exists  $v \in \mathcal{C}^{\text{box}}(R)$  such that  $q \in \text{bc}_R(v)$ , hence  $p$  and  $p'$  are the unique auxiliary port of  $v$ , therefore  $p = p'$ .
2. By induction on the length  $n \in \mathbb{N}$  of the box-crossing path  $(p_i)_{i \in I}$  from  $p$  to  $q$ . If  $n = 0$  then  $p = p_0 = q$ , thus there exists a box-crossing path from  $p'$  to  $p$  by hypothesis. If  $n > 0$  then there are only two cases: if  $p' = q$  then there exists a box-crossing path from  $p$  to  $p'$  by hypothesis; otherwise there exist  $q' \in \mathcal{P}(\Phi)$ , a box-crossing path from  $p'$  to  $q'$  and a box-crossing path of length 1 from  $q'$  to  $q$ , so  $p_{n-1} = q'$  by lemma 49.1, therefore there exists a box-crossing path from  $p$  to  $p'$  or from  $p'$  to  $p$  by induction hypothesis applied to  $p_{n-1}$ .
3. By definition 30, there is no ascending path of length 1 from  $q$  to  $c$ , hence the only possibility is that the path from  $q$  to  $c$  has length 0, therefore  $q = c$ .
4. Let us suppose that  $(p_i)_{0 \leq i \leq m}$  and  $(q_j)_{0 \leq j \leq n}$  (for some  $m, n \in \mathbb{N}$ ) are two ascending paths such that  $p_0 = p = q_0$  and  $p_m = q = q_n$ : we prove by induction on  $m$  that  $m = n$  and  $p_i = q_i$  for every  $0 \leq i \leq m$ .
 

If  $m = 0$  then  $p = q_0 = p_0 = q$ , furthermore  $n = 0$  (otherwise  $q = q_0 \preceq_R^1 q_1$  and  $q_1 \preceq_R q_n = q$  with  $q \neq q_1$ , that is impossible since  $\preceq_R$  is antisymmetric by proposition 47.1).

If  $m > 0$  then  $p_{m-1} = q_{n-1}$  by lemma 49.1. By induction hypothesis,  $m - 1 = n - 1$  and  $p_i = q_i$  for every  $0 \leq i \leq m - 1$ , thus we can conclude.
5. For the unicity, if  $c, c' \in \mathcal{P}^{\text{free}}(R) \cup \bigcup \mathcal{Cuts}_0(R) \cup \mathcal{D}_0(R)$  are such that there exist two box-crossing path in  $R$  from  $c$  to  $q$  and from  $c'$  to  $q$  then there exists a box-crossing path in  $R$  from  $c$  to  $c'$  or from  $c'$  to  $c$  by lemma 49.2, hence  $c = c'$  by lemma 49.3.

For the existence, we build an initial segment  $J$  of  $\mathbb{N}$  and a sequence of finite paths  $(\varphi_j)_{j \in J}$  in  $\Phi$  as follows:



- $0 \in J$  and  $\varphi_0$  is the path of length 0 consisting only of  $q$ ;
- if  $j \in J$  and  $\varphi_j = (p_i)_{0 \leq i \leq n}$  then:
  - if  $p_0 \in \mathcal{P}^{\text{free}}(\Phi) \cup \bigcup \mathcal{Cuts}_0(\Phi) \cup \mathcal{D}_0(R)$  then  $J = \{0, \dots, j\}$ ,
  - otherwise,  $j+1 \in J$  and  $\varphi_{j+1} = (q_i)_{0 \leq i \leq n+1}$  where  $q_{i+1} = p_i$  for every  $0 \leq i \leq n$  and
    - \* if  $p_0 \in \mathcal{P}_l^{\text{aux}}(\Phi)$  for some  $l \in \mathcal{C}(\Phi)$  then  $q_0 = \mathbf{P}_\Phi^{\text{pri}}(l)$ ;
    - \* if  $p_0 \in \mathcal{P}^{\text{pri}}(\Phi)$  then either  $q_0 \in \mathcal{P}^{\text{aux}}(\Phi)$  with  $\{p_0, q_0\} \in \mathcal{W}(\Phi)$ , or  $q_0$  is the unique auxiliary port of some  $v \in \mathcal{C}^{\text{box}}(\Phi)$  such that  $p_0 \in \text{cutports}_\Phi(v)$ ;
    - \* if  $p_0 \in \mathcal{D}(R)$  with  $p_0 \in \text{bc}_R(v)$  for some  $v \in \mathcal{C}^{\text{box}}(R)$  then  $q_0$  is the unique auxiliary port of  $v$ .

By construction, for every  $j \in J$ ,  $\varphi_j$  is a box-crossing path in  $R$  of length  $j$  ending in  $q$  and moreover, if  $j+1 \in J$  then  $\varphi_j$  is a sub-path of  $\varphi_{j+1}$ .  $J$  is a finite set (otherwise there would be an infinite box-crossing path  $\varphi_\omega$  such that, for every  $j \in J = \mathbb{N}$ ,  $\varphi_j$  would be a sub-path of  $\varphi_\omega$ , that is impossible since  $\mathcal{P}(\Phi)$  is a finite set and  $\preceq_R$  is antisymmetric), hence there exists  $m \in \mathbb{N}$  such that  $J = \{0, \dots, m\}$  and  $\varphi_m$  is a box-crossing path from  $c \in \mathcal{P}^{\text{free}}(R) \cup \bigcup \mathcal{Cuts}_0(R) \cup \mathcal{D}_0(R)$  to  $q$ .  $\square$

Notice that lemmas 49.1,2,4,5 do not hold in the case of generic ascending paths (i.e. for the order relation  $\preceq_R$ ).

**Definition 50.** Let  $R \in \mathbf{PS}$  and  $p \in \mathcal{P}(R)$ .

We denote by  $\text{c}_R(p)$  the unique  $c \in \mathcal{P}^{\text{free}}(R) \cup \bigcup \mathcal{Cuts}_0(R) \cup \mathcal{D}_0(R)$  such that  $c \preceq_R p$ .

We set  $\text{boxesof}_R(p) = \{v \in \mathcal{C}^{\text{box}}(R) \mid p \in \text{inbox}_R(v)\}$ . If  $\text{boxesof}_R(p) \neq \emptyset$ , we denote by  $\mathbf{C}_R^{\text{box}}(p)$  the  $v \in \text{boxesof}_R(p)$  such that  $\text{inbox}_R(v)$  is minimal with respect to  $\subseteq$ .

The ground of  $R$  is  $\text{Ground}(R) = \{q \in \mathcal{P}(R) \mid \text{boxesof}_R(q) = \emptyset\}$ .

By lemma 49.5, the function  $\text{c}_R$  is well-defined for any  $R \in \mathbf{PS}$ . By the nesting condition,  $\mathbf{C}_R^{\text{box}}(p)$  is well-defined for every  $R \in \mathbf{PS}$  and  $p \in \mathcal{P}(R)$  such that  $p \in \text{inbox}_R(v)$  for some  $v \in \mathcal{C}^{\text{box}}(R)$ .

Given  $R \in \mathbf{PS}$  and  $p \in \mathcal{P}(R)$ ,  $\text{boxesof}_R(p)$  is morally the set of boxes in  $R$  containing  $p$ .

In spite of our non-inductive definition of proof-structure, we can recover some typical informations of the inductive one, such as the depth of a port and the depth of a proof-structure. The following definition is nothing but the adaptation to our syntax (allowing ps with cuts and deadlocks) of the corresponding definition in [dCT12].

**Definition 51 (Depth).** Let  $R \in \mathbf{PS}$ .

Let  $p \in \mathcal{P}(R)$  and let  $\varphi_p$  be the box-crossing path from  $\mathbf{c}_R(p)$  to  $p$ . The depth of  $p$  in  $R$ , denoted by  $\text{depth}_R(p)$ , is a nonnegative integer defined by:

$$\text{depth}_R(p) = \text{card}(\{l \in \mathcal{C}^{\text{box}}(R) \mid \varphi_p \text{ crosses } l\}) + \sum_{q \in \text{Auxdoors}(R)} \text{card}(\{l' \in \mathcal{C}^{\text{box}}(R) \mid \varphi_p \text{ crosses } q \in \text{auxd}_R(l')\}).$$

If  $\text{depth}_R(p) = n$  then we say that “ $p$  is at depth  $n$  in  $R$ ”.

For every  $l \in \mathcal{C}(R)$ , the depth of  $l$  in  $R$ , denoted by  $\text{depth}_R(l)$ , is the depth of  $\mathbf{P}_R^{\text{pri}}(l)$  in  $R$ . If  $\text{depth}_R(l) = n$  then we say that “ $l$  is at depth  $n$  in  $R$ ”.

For every  $v \in \mathcal{C}^{\text{box}}(R)$ , the depth of  $\text{box}_R(v)$  in  $R$ , denoted by  $\text{depth}_R(\text{box}_R(v))$ , is the depth of  $v$  in  $R$ . If  $\text{depth}_R(\text{box}_R(v)) = n$  then we say that “ $\text{box}_R(v)$  is at depth  $n$  in  $R$ ”.

The depth of  $R$  is  $\text{depth}(R) = \sup\{\text{depth}_R(p) \mid p \in \mathcal{P}(R)\}$ .

The depth of a port  $p$  in a ps  $R$  is well-defined thanks to lemma 49.4, which says that there exists a unique box-crossing path  $\varphi_p$  from  $\mathbf{c}_R(p)$  to  $p$ . Roughly speaking,  $\text{depth}_R(p)$  is calculated by counting the number of promotion cells and arrows pointing to auxiliary doors of  $R$  crossed by  $\varphi_p$ .

**Remark 52.** Let  $R \in \mathbf{PS}$ .

1. If  $R \in \mathbf{PPS}_{\text{DiLL}_0}$  then  $\text{depth}(R) = 0$ , as  $\mathcal{C}^{\text{box}}(R) = \emptyset$  (remember that  $\mathbf{PPS}_{\text{DiLL}_0} \subseteq \mathbf{PS}$ ).
2. Let  $v \in \mathcal{C}^{\text{box}}(R)$ : if  $p_v$  is the unique auxiliary port of  $v$ , then  $\text{depth}_R(p_v) = \text{depth}_R(\text{box}_R(v)) + 1$ . Indeed if  $\varphi_{\mathbf{P}_R^{\text{pri}}(v)}$  is the unique box-crossing path from  $\mathbf{c}_R(\mathbf{P}_R^{\text{pri}}(v)) = \mathbf{c}_R(p_v)$  to  $\mathbf{P}_R^{\text{pri}}(v)$ , then  $\varphi_{\mathbf{P}_R^{\text{pri}}(v)} \cdot p_v$  is the unique box-crossing path from  $\mathbf{c}_R(p_v)$  to  $p_v$  and moreover  $\varphi_{\mathbf{P}_R^{\text{pri}}(v)} \cdot p_v$  crosses  $v$  whereas  $\varphi_{\mathbf{P}_R^{\text{pri}}(v)}$  does not.
3. Let  $v \in \mathcal{C}^{\text{box}}(R)$ : if  $p \in \text{inbox}_R(v) \cap \mathcal{P}(\text{box}_R(v))$  then  $\text{boxesof}_{\text{box}_R(v)}(p) = \text{boxesof}_R(p) \setminus \{l \in \mathcal{C}^{\text{box}}(R) \mid \mathbf{P}_R^{\text{pri}}(l) \notin \text{inbox}_R(v)\}$ . Indeed, let  $l \in \mathcal{C}^{\text{box}}(R)$ :  $l \in \text{boxesof}_{\text{box}_R(v)}(p)$  iff  $l \in \mathcal{C}^{\text{box}}(\text{box}_R(v))$  and  $p \in \text{inbox}_{\text{box}_R(v)}(l) = \text{inbox}_R(l)$  iff  $\mathbf{P}_R^{\text{pri}}(l) \in \text{inbox}_R(v)$  and  $l \in \text{boxesof}_R(v)$ .

The following lemma shows that the depth of a port  $p$  in a ps  $R$  is nothing but the number of boxes in  $R$  containing  $p$ , as in the usual inductive syntaxes of linear logic.

**Lemma 53.** Let  $R \in \mathbf{PS}$ . For every  $p \in \mathcal{P}(R)$ , one has

$$\text{depth}_R(p) = \text{card}(\text{boxesof}_R(p))$$

In particular,  $\text{depth}_R(p) = 0$  iff  $p \notin \text{inbox}_R(v)$  for any  $v \in \mathcal{C}^{\text{box}}(R)$ .

PROOF. By induction on the  $\text{depth}(p) \in \mathbb{N}$ . We denote by  $\varphi_p = (p_i)_{0 \leq i \leq n}$  (with  $n \in \mathbb{N}$ ) the unique box-crossing path from  $c_R(p)$  to  $p$ .

If  $\text{depth}_R(p) = 0$ , then  $\varphi_p$  does not cross any  $v \in \mathcal{C}^{\text{prom}}(R) = \mathcal{C}^{\text{box}}(R)$  (as  $R \in \mathbf{PS}$ ) nor any  $q \in \text{Auxdoors}(R)$ , according to definition 51. Hence, for every  $v \in \mathcal{C}^{\text{box}}(R)$ , every path above  $v$  does not end in  $p$ , by condition 4 in definition 38. Thus,  $\text{card}(\text{boxesof}_R(p)) = 0 = \text{depth}_R(p)$ .

If  $\text{depth}_R(p) > 0$ , then  $\varphi_p$  crosses some  $v \in \mathcal{C}^{\text{box}}(R)$  or some  $q \in \text{Auxdoors}(R)$ , according to definition 51. Hence,  $n > 0$  and for some  $v \in \mathcal{C}^{\text{box}}(R)$ , there exists a path above  $v$  ending in  $p$  (i.e.  $p \in \text{inbox}_R(v)$ ), by condition 4 in definition 38. According to definition 50 (i.e. thanks to nesting condition, as  $R \in \mathbf{PS}$ ),  $\mathcal{C}_R^{\text{box}}(p)$  is defined. By condition 4 in definition 38 and lemma 49.4,  $\varphi_p$  crosses either  $\mathcal{C}_R^{\text{box}}(p)$  or  $q \in \text{auxd}_R(\mathcal{C}_R^{\text{box}}(v))$ : in both cases, there exists  $k \in \{0, \dots, n\}$  such that  $p_k \in \text{doors}_R(\mathcal{C}_R^{\text{box}}(p))$  and the subpath  $(p_i)_{k \leq i \leq n}$  of  $\varphi_p$  does not cross any  $v \in \mathcal{C}^{\text{box}}(R)$  nor any  $q \in \text{Auxdoors}(R)$ , because of the minimality of  $\text{inbox}_R(\mathcal{C}_R^{\text{box}}(p))$ . So  $\text{depth}_R(p_{k-1}) = \text{depth}_R(p) - 1$ , hence  $\text{depth}_R(p_{k-1}) = \text{card}(\text{boxesof}_R(p_{k-1}))$  by induction hypothesis. Because of lemma 43,  $p_{k-1} \notin \text{inbox}_R(\mathcal{C}_R^{\text{box}}(p))$  whereas  $p_k \in \text{inbox}_R(\mathcal{C}_R^{\text{box}}(p))$  and thus  $\text{card}(\text{boxesof}_R(p_{k-1})) = \text{card}(\text{boxesof}_R(p)) - 1$ , by the nesting condition. Therefore,  $\text{depth}_R(p) = \text{card}(\text{boxesof}_R(p))$ .  $\square$

The following proposition reveals some intuitive properties of some notions already introduced.

**Proposition 54.** *Let  $R \in \mathbf{PS}$ .*

1. *If  $\{p, q\} \in \mathcal{W}(R)$  then  $\text{boxesof}_R(p) = \text{boxesof}_R(q)$  and  $\text{depth}_R(p) = \text{depth}_R(q)$ .*
2. *Let  $v \in \mathcal{C}^{\text{box}}(R)$  and let  $p_v$  be the unique auxiliary port of  $v$ : for every  $q \in \text{bc}_R(v)$  one has  $\text{boxesof}_R(p_v) = \text{boxesof}_R(q)$  and  $\text{depth}_R(q) = \text{depth}_R(p_v)$ .*
3. *Let  $p \in \bigcup \text{Cuts}(R) \cup \mathcal{D}(R)$ :  $\text{depth}_R(p) = 0$  iff  $p \in \bigcup \text{Cuts}_0(R) \cup \mathcal{D}_0(R)$ .*
4. *Let  $v, v' \in \mathcal{C}^{\text{box}}(R)$ : if  $\text{inbox}_R(v') \subseteq \text{inbox}_R(v)$ , then  $\text{boxesof}_R(\text{P}_R^{\text{pri}}(v)) \subseteq \text{boxesof}_R(\text{P}_R^{\text{pri}}(v'))$  and  $\text{depth}_R(\text{box}_R(v)) \leq \text{depth}_R(\text{box}_R(v'))$ .*
5. *for every  $v \in \mathcal{C}^{\text{box}}(R)$  one has  $\text{depth}(\text{box}_R(v)) < \text{depth}(R)$ .*

PROOF.

1. For every  $v \in \mathcal{C}^{\text{box}}(R)$ , there exists a path above  $v$  ending in  $p$  iff there exists a path above  $v$  ending in  $q$ : this is evident when  $\{p, q\}$  is neither a cut nor an axiom, this is due to condition 4 in definition 38 if  $\{p, q\}$  is an axiom, and this is due to definition of  $\text{bc}_R$  if  $\{p, q\}$  is a cut. Hence,  $p \in \text{inbox}_R(v)$  iff  $q \in \text{inbox}_R(v)$  and thus  $\text{boxesof}_R(p) = \text{boxesof}_R(q)$ . Therefore  $\text{depth}_R(p) = \text{depth}_R(q)$  by lemma 53.

2. For every  $l \in \mathcal{C}^{\text{box}}(R)$ ,  $l \in \text{boxesof}_R(p_v)$  iff  $p_v \in \text{inbox}_R(l)$  iff there exists in  $R$  a path above  $l$  ending in  $p_v$  iff (by definition of ascending path) there exists in  $R$  a path above  $l$  ending in  $q$  iff  $q \in \text{inbox}_R(l)$  iff  $l \in \text{boxesof}_R(q)$ . Therefore,  $\text{boxesof}_R(p_v) = \text{boxesof}_R(q)$  and so  $\text{depth}_R(q) = \text{depth}_R(p_v)$  by lemma 53.
3. If  $\text{depth}_R(p) = 0$  then the box-crossing path from  $\text{c}_R(p)$  to  $p$  crosses no promotion cells, hence there is no ports  $q$  such that  $q \preceq_R^1 p$ , in particular  $p \notin \text{im}(\text{bc}_R)$  and thus  $p \in \text{Cuts}_0(R) \cup \mathcal{D}_0(R)$ .  
Conversely, if  $p \in \text{Cuts}_0(R) \cup \mathcal{D}_0(R)$  then  $p = \text{c}_R(p)$  by lemma 49.3, therefore  $\text{depth}_R(p) = 0$ .
4. By remark 42.2, in  $R$  there exists a path above  $v$  ending in the unique auxiliary port of  $v'$ . If  $\text{P}_R^{\text{pri}}(v') \notin \text{inbox}_R(v)$  then there is no path above  $v$  ending in  $\text{P}_R^{\text{pri}}(v')$ , hence  $v = v'$  and so  $\text{P}_R^{\text{pri}}(v) = \text{P}_R^{\text{pri}}(v')$ , whence  $\text{boxesof}_R(\text{P}_R^{\text{pri}}(v)) = \text{boxesof}_R(\text{P}_R^{\text{pri}}(v'))$ . Otherwise in  $R$  there exists a path  $\varphi$  above  $v$  ending in  $\text{P}_R^{\text{pri}}(v')$ : if  $l \in \text{boxesof}_R(\text{P}_R^{\text{pri}}(v))$  then  $l \in \mathcal{C}^{\text{box}}(R)$  and  $\text{P}_R^{\text{pri}}(v) \in \text{inbox}_R(l)$ , thus there exists a path  $\psi$  above  $l$  ending in  $\text{P}_R^{\text{pri}}(v)$  and so  $\psi \cdot \varphi$  is a path in  $R$  above  $l$  ending in  $\text{P}_R^{\text{pri}}(v')$ , whence  $\text{P}_R^{\text{pri}}(v') \in \text{inbox}_R(l)$  and thus  $l \in \text{boxesof}_R(\text{P}_R^{\text{pri}}(v'))$ . Therefore  $\text{boxesof}_R(\text{P}_R^{\text{pri}}(v)) = \text{boxesof}_R(\text{P}_R^{\text{pri}}(v'))$ , so  $\text{depth}_R(\text{box}_R(v)) \leq \text{depth}_R(\text{box}_R(v'))$  by lemma 53.
5. Let  $p \in \mathcal{P}(\text{box}_R(v))$  be such that  $\text{depth}(\text{box}_R(v)) = \text{depth}_{\text{box}_R(v)}(p)$ . If  $p \in \mathcal{P}(\text{box}_R(v)) \setminus \text{inbox}_R(v)$  then  $p$  is the principal port of a unary ?-cell whose unique auxiliary port  $q \in \text{auxd}_R(v)$  and so  $q \in \text{inbox}_R(v)$  and  $\text{depth}_{\text{box}_R(v)}(p) \leq \text{depth}_{\text{box}_R(v)}(q)$ . Therefore, we can suppose without loss of generality that  $p \in \text{inbox}_R(v)$  and thus  $\text{boxesof}_{\text{box}_R(v)}(p) = \text{boxesof}_R(p) \setminus \{l \in \mathcal{C}^{\text{box}}(R) \mid \text{P}_R^{\text{pri}}(l) \notin \text{inbox}_R(v)\}$  by remark 52.3. As  $\text{P}_R^{\text{pri}}(v) \notin \text{inbox}_R(v)$  (by lemma 43), one has  $\text{boxesof}_{\text{box}_R(v)}(p) \subsetneq \text{boxesof}_R(p)$ , whence  $\text{depth}(\text{box}_R(v)) = \text{depth}_{\text{box}_R(v)}(p) < \text{depth}_R(p) \leq \text{depth}(R)$  by lemma 53.  $\square$

According to proposition 54.1, we are entitled to talk about of the depth of a wire  $\{p, q\}$  of a ps: it is the depth of  $p$  or  $q$ . Propositions 54.2-3 mean that in a ps  $R$ , for any  $v \in \mathcal{C}^{\text{box}}(R)$ , the cuts and deadlocks pointed by the arrow function  $\text{bc}_R$  are the cuts and deadlock at depth 0 in  $\text{box}_R(v)$ , in other words  $\text{box}_R(v)$  is the deepest (i.e. “smallest” in the sense of proposition 54.4) box containing them.

Notice that proposition 54.2 does not hold if we replace the hypothesis  $q \in \text{bc}_R(v)$  with  $q \in \text{auxd}_R(v)$  because in general, we can have  $v, l \in \mathcal{C}^{\text{box}}(R)$  such that  $q \in \text{auxd}_R(l) \cap \text{auxd}_R(v)$  but  $p_v \notin \text{inbox}_R(l)$  (where  $p_v$  is the unique auxiliary port of  $v$ ), whence  $\text{boxesof}_R(q) \not\subseteq \text{boxesof}_R(p_v)$ .

Proposition 54.5 allows to make easily induction on the depth of a ps.

## 1.6 Indexed ((pre-)pre-)proof-structures

We introduce the notion of indexed pseudo-structure (resp. ppps; pps; ps), i.e. a pseudo-structure (resp. ppps; pps; ps) with ordered conclusions. This is mandatory to fix an order on conclusions in order to define the interpretation of a proof-structure in the relational model.

**Definition 55** (Indexed ((pre-)pre-)proof-structure). *An indexed pseudo-structure is a pair  $(\Phi, \text{ind})$  such that  $\Phi \in \mathbf{PseudoPPPS}$  and  $\text{ind} : \mathcal{P}^{\text{free}}(\Phi) \rightarrow \{1, \dots, \text{card}(\mathcal{P}^{\text{free}}(\Phi))\}$  is a bijection. We say then that  $\text{ind}$  is an enumeration of  $\mathcal{P}^{\text{free}}(\Phi)$ .*

*An indexed ppps (resp. indexed pps; indexed ps) is an indexed pseudo-structure  $(\Phi, \text{ind})$  such that  $\Phi \in \mathbf{PPPS}$  (resp.  $\Phi \in \mathbf{PPS}$ ;  $\Phi \in \mathbf{PS}$ ).*

*We denote by  $\mathbf{PseudoPPPS}^{\text{ind}}$  (resp.  $\mathbf{PPPS}^{\text{ind}}$ ;  $\mathbf{PPS}^{\text{ind}}$ ;  $\mathbf{PS}^{\text{ind}}$ ) the set of indexed pseudo-structures (resp. indexed ppps; indexed pps; indexed ps).*

*We set  $\mathbf{PPS}_{\text{DILL}_0}^{\text{ind}} = \{(R, \text{ind}) \in \mathbf{PS}^{\text{ind}} \mid R \in \mathbf{PPS}_{\text{DILL}_0}\}$  and  $\mathbf{PS}_{\text{MELL}}^{\text{ind}} = \{(R, \text{ind}) \in \mathbf{PS}^{\text{ind}} \mid R \in \mathbf{PS}_{\text{MELL}}\}$ .*

We introduce the notion of “identity” (or better said isomorphism) between two indexed ppps (resp. pps; ps). The idea is that two corresponding conclusions of two indexed ppps (resp. pps; ps) have to be in the same order position.

**Definition 56** (Isomorphism between indexed ((pre-)pre-)proof-structures). *Let  $(\Phi, \text{ind}), (\Phi', \text{ind}') \in \mathbf{PPPS}^{\text{ind}}$  (resp.  $(\Phi, \text{ind}), (\Phi', \text{ind}') \in \mathbf{PPS}^{\text{ind}}$ ;  $(\Phi, \text{ind}), (\Phi', \text{ind}') \in \mathbf{PS}^{\text{ind}}$ ).*

*An isomorphism from  $(\Phi, \text{ind})$  to  $(\Phi', \text{ind}')$  is a  $\varphi : \Phi \simeq \Phi'$  such that the following diagram commutes:*

$$\begin{array}{ccc}
 \mathcal{P}^{\text{free}}(\Phi) & \xrightarrow{\text{ind}} & \{1, \dots, \text{card}(\mathcal{P}^{\text{free}}(\Phi))\} \\
 \varphi_{\mathcal{P}} \downarrow & & \nearrow \text{ind}' \\
 \mathcal{P}^{\text{free}}(\Phi') & & 
 \end{array}$$

*We write then  $\varphi : (\Phi, \text{ind}) \simeq (\Phi', \text{ind}')$ .*

*If there exists an isomorphism from  $(\Phi, \text{ind})$  to  $(\Phi', \text{ind}')$ , then we say that  $(\Phi, \text{ind})$  and  $(\Phi', \text{ind}')$  are isomorphic and we write  $(\Phi, \text{ind}) \simeq (\Phi', \text{ind}')$ .*

**Remark 57.** Let  $\Phi, \Phi' \in \mathbf{PPPS}$  with  $\varphi : R \simeq R'$ . For every enumeration  $\text{ind}$  of  $\mathcal{P}^{\text{free}}(\Phi)$ , there exists an enumeration  $\text{ind}'$  of  $\mathcal{P}^{\text{free}}(\Phi')$  such that  $\varphi : (\Phi, \text{ind}) \simeq (\Phi', \text{ind}')$ . Indeed, it suffices to take  $\text{ind}' : \mathcal{P}^{\text{free}}(\Phi') \rightarrow \{1, \dots, \text{card}(\mathcal{P}^{\text{free}}(\Phi'))\}$  such that  $\text{ind}'(p) = \text{ind}(\varphi_{\mathcal{P}}^{-1}(p))$ .

## 1.7 A non-inductive correctness criterion

A correctness criterion is a property fulfilled by all and only those proof-structures corresponding to a proof in the (multiplicative and exponential framework of) Linear Logic sequent calculus. This gives a geometrical account of the Linear Logic proofs. There is a multitude of equivalent correctness criteria for the multiplicative and exponential framework of Linear Logic, the most common one is the Danos-Regnier criterion (see for example [DR89, Tor03]), which is a simplification of the primary long trip criterion of Girard introduced in [Gir87].

All the well-known correctness criteria for the multiplicative and exponential framework of Linear Logic proof-structures are defined by induction on the depth of the proof-structure, so they can be considered “purely geometrical” only in the case of a proof-structure without boxes (in particular, in the multiplicative framework). In our syntax we can reformulate the Danos-Regnier correctness criterion for the multiplicative and exponential framework of Linear Logic proof-structures in such a way that our criterion is completely “non-inductive”, that is reinforces the idea of a “purely geometrical” characterization of (multiplicative and exponential) Linear Logic proofs.

**Definition 58** (Snipping, linearization). *Let  $\Phi \in \mathbf{Modules}$  and let  $Q \subseteq \mathcal{P}^{\text{aux}}(\Phi)$ . The snipping of  $Q$  in  $\Phi$  is a 6-tuple  $\Phi' = (\mathcal{C}', \mathcal{I}', \mathcal{D}', \mathcal{W}', \text{auxd}', \text{bc}')$  such that:*

- $\mathcal{C}' = (\mathbf{t}_\Phi, \mathcal{P}', \mathbf{C}_\Phi \upharpoonright_{\mathcal{P}'}, \mathcal{P}_\Phi^{\text{pri}}, \mathcal{P}_\Phi^{\text{left}} \upharpoonright_{L_Q})$  where  $\mathcal{P}' = \mathcal{P}(\Phi) \setminus Q$  and  $L_Q = \{l \in \mathcal{C}^{\otimes, \exists}(\Phi) \mid \text{card}(\{p \in \mathcal{P}' \mid \mathbf{C}_\Phi(p) = l\}) = 3\}$ ;
- $\mathcal{I}' = \mathcal{I}(\Phi) \setminus (Q \cap \bigcup \mathcal{A}x(\Phi))$ ;
- $\mathcal{D}' = \mathcal{D}(\Phi)$  and  $\mathcal{W}' = \mathcal{W}(\Phi)$ ;
- $\text{auxd}' = \emptyset = \text{bc}'$  (where  $\emptyset$  is the empty function).

If  $Q = \emptyset$ , we say that the snipping of  $Q$  in  $\Phi$  is the linearization of  $\Phi$ , denoted by  $\text{lin}(\Phi)$ .

**Remark 59.** For every  $\Phi \in \mathbf{Modules}$  and  $Q \subseteq \mathcal{P}^{\text{aux}}(\Phi)$ , if  $\Phi' = (\mathcal{C}', \mathcal{I}', \mathcal{D}', \mathcal{W}', \text{auxd}', \text{bc}')$  is the snipping of  $Q$  in  $\Phi$ , then  $\mathcal{C}' \in \mathbf{ModuleBases}$  and  $\Phi' \in \mathbf{Module}$ . Furthermore  $\mathcal{C}(\Phi) = \mathcal{C}(\Phi')$ ,  $\mathcal{P}(\Phi) = \mathcal{P}(\Phi')$ ,  $\mathcal{P}^{\text{pri}}(\Phi) = \mathcal{P}^{\text{pri}}(\Phi')$ ,  $\mathcal{D}(\Phi) = \mathcal{D}_0(\Phi')$  and  $\mathcal{W}(\Phi) = \mathcal{W}(\Phi')$ , but in general  $\mathcal{I}(\Phi) \not\subseteq \mathcal{I}(\Phi')$ ,  $\mathcal{P}^{\text{free}}(\Phi) \not\subseteq \mathcal{P}^{\text{free}}(\Phi')$  and  $\mathcal{P}^{\text{aux}}(\Phi) \not\subseteq \mathcal{P}^{\text{aux}}(\Phi')$ . Therefore, every path in  $\Phi'$  is also a path in  $\Phi$  (but the converse does not hold).

If  $R \in \mathbf{PS}$ , then  $\text{lin}(R) \in \mathbf{PPS}_{\text{DiLL}_0}$ ,  $\mathcal{I}(R) = \mathcal{I}(\text{lin}(R))$ ,  $\mathcal{P}^{\text{free}}(R) = \mathcal{P}^{\text{free}}(\text{lin}(R))$  and  $\mathcal{P}^{\text{aux}}(R) = \mathcal{P}^{\text{aux}}(\text{lin}(R))$ .

Roughly speaking, if  $\Phi$  is a module and  $Q$  is a set of auxiliary ports of  $\Phi$ , the snipping of  $Q$  in  $\Phi$  is the module obtained from  $\Phi$  by disconnecting

the ports of  $Q$  (which become isolated ports) by their cells. This operation might transform a ps in a module which is not a pseudo-structure, because some binary  $\otimes$ - or  $\mathfrak{Y}$ -cells might get 0-ary or unary, or some wires might get hanging (i.e. they connect a principal port with an isolated port).

The linearization of a ps  $\Phi$  has to be seen as the  $\text{DiLL}_0$ -proof-structure obtained from  $\Phi$  by forgetting the boundaries of all the boxes of  $\Phi$  (i.e. their arrows).

**Definition 60** (Switching, correctness graph). *Let  $R \in \mathbf{PS}$  and let*

$$\mathcal{C}^{\mathfrak{Y},?c}(R) = \mathcal{C}^{\mathfrak{Y}}(R) \cup \{l \in \mathcal{C}(R) \mid \mathfrak{t}_R(l) = ?, \mathfrak{a}_R(l) \geq 2\}.$$

*A switching of  $R$  is a function associating with every  $l \in \mathcal{C}^{\mathfrak{Y},?c}(R)$  an auxiliary port of  $l$ .*

*For every switching  $s$  of  $R$ , we set  $\text{off}_R(s) = \{p \in \mathcal{P}_l^{\text{aux}}(R) \setminus \text{im}(s) \mid l \in \mathcal{C}^{\mathfrak{Y},?c}(R)\}$ : the  $s$ -correctness graph of  $\Phi$  is the snipping of  $\text{off}_R(s)$  in  $R$ .*

*A correctness graph of  $R$  is a  $s$ -correctness graph of  $R$  for some switching  $s$  of  $R$ .*

Definition 60 reformulates in our syntax some standard notions of Linear Logic proof-structures.

**Remark 61.** For every  $R \in \mathbf{PS}$  and switching  $s$  of  $R$ , if  $G_R^s$  is the  $s$ -correction graph of  $R$  then  $G_R^s \in \mathbf{Modules}$ ; moreover, if  $l \in \mathcal{C}(R)$  is such that there exists a  $p \in \text{doors}_R(v) \cap \mathcal{P}_l^{\text{aux}}(R)$  for some  $v \in \mathcal{C}^{\text{box}}(R)$ , then  $\mathfrak{a}_{G_R^s}(l) = 1$ .

**Conjecture 62** (Cryptic). Let  $R \in \mathbf{PS}$ .  $R$  is acyclic iff for every switching  $s$  in  $\text{boxed}(R)$ ,  $l \in \mathcal{C}^{\text{box}}(\text{boxed}(R))$ ,  $p \in \mathcal{P}(\text{boxed}(R))$  there exists at most one path in  $s$  above  $l$  ending in  $p$ .

WHY?

**Definition 63** (DR-path). *Let  $R \in \mathbf{PS}$  and let  $s$  be a switching of  $R$ .*

*A DR-path in  $R$  according to  $s$  is a path  $(p_i)_{i \in I}$  (where  $I$  is an initial segment of  $\mathbb{N}$ ) in the  $s$ -correction graph of  $R$  such that for every  $v \in \mathcal{C}^{\text{box}}(R)$  and  $i, i+1 \in I$ , if  $p_i \in \text{doors}_R(v) \cap \mathcal{P}_l^{\text{aux}}(R)$  and  $p_{i+1} = \mathbf{P}_R^{\text{pr}_i}(l)$  for some  $l \in \mathcal{C}(R)$ , then for every  $j \in I$  such that  $j > i+1$  one has  $p_j \notin \text{doors}_R(v)$ .*

*For every  $p, q \in \mathcal{P}(R)$ , we say that  $p$  and  $q$  are DR-connected according to  $s$  if there exists a DR-path in  $R$  according to  $s$  from  $p$  to  $q$ .*

The idea is that, given a ps  $R$  and a switching  $s$  of  $R$ , a DR-path in  $R$  according to  $s$  leaving a box cannot re-enter it. By means of DR-paths we can give a simple correctness criterion.

**Definition 64** (DR-connection, DR-acyclicity, proof-net). *Let  $R \in \mathbf{PS}$ .*

*$R$  is DR-connected if for every switching  $s$  of  $R$ , all  $p, q \in \mathcal{P}(R)$  are DR-connected according to  $s$ .*

*$R$  is DR-acyclic if  $\mathcal{D}(R) = \emptyset$  and, for every switching  $s$  of  $R$ , each path in the  $s$ -correction graph of  $R$  is not a cycle and it is a DR-path in  $R$  according to  $s$ .*



$R$  is ACC (or  $R$  is a proof-net or  $R$  satisfies the correctness criterion) if  $R$  is DR-connected and DR-acyclic.

We denote by  $\mathbf{PN}$  the set of proof-nets. We set  $\mathbf{PN}_{\text{MELL}} = \mathbf{PN} \cap \mathbf{PS}_{\text{MELL}}$  (resp.  $\mathbf{PN}_{\text{DiLL}_0} = \mathbf{PN} \cap \mathbf{PPS}_{\text{DiLL}_0}$ ), whose elements are the MELL-proof-nets (resp. DiLL<sub>0</sub>-proof-nets).

We point out that our correctness criterion, as well as our definition of proof-structure, is completely non-inductive, by taking into account only some “geometrical informations” available in a proof-structure. Our correctness criterion can be seen as a non-inductive version of the Danos-Regnier correctness criterion [DR89, Tor03] for the multiplicative and exponential framework of Linear Logic (which is non-inductive only in the multiplicative fragment). Intuitively, in our correctness criterion DR-paths play the role of induction on the depth of a proof-structure in the Danos-Regnier criterion.

In the case of a DR-connected ps, we can simplify the correctness criterion.

**Proposition 65.** *Let  $R \in \mathbf{PS}$  be DR-connected.  $R$  is a proof-net iff  $\mathcal{D}(R) = \emptyset$  and, for every switching  $s$  of  $R$ , the  $s$ -correction graph of  $R$  is acyclic.*

PROOF. Let  $R \in \mathbf{PS}$  be DR-connected such that  $\mathcal{D}(R) = \emptyset$  and, for every switching  $s$  of  $R$ , the  $s$ -correctness graph of  $R$  is acyclic. We have to show that each path in the  $s$ -correction graph of  $R$  is a DR-path in  $R$  according to  $s$ . Let us suppose by absurd that for some switching  $s$  of  $R$  there exists a path in the  $s$ -correction graph  $G_R^s$  of  $R$  which is not a DR-path in  $R$  according to  $s$ . Thus, there would exist  $v \in \mathcal{C}^{\text{box}}(R)$  and a path  $(p_i)_{0 \leq i \leq n}$  (for some  $n \in \mathbb{N}^*$ ) in  $G_R^s$  the such that  $p_n \in \text{doors}_R(v)$ ,  $p_0 \in \text{doors}_R(v) \cap \mathcal{P}_l^{\text{aux}}(R)$  and  $p_1 = \text{P}_R^{\text{pri}}(l)$  for some  $l \in \mathcal{C}(R)$  and  $p_i \notin \text{doors}_R(v)$  for every  $1 \leq i \leq n-1$ . Since  $R$  is DR-connected, there would exist a DR-path  $(q_j)_{0 \leq j \leq m}$  (for some  $m \in \mathbb{N}^*$ ) in  $R$  according to  $s$  from  $p_n$  to  $p_0$ , hence  $q_1 \in \text{inbox}_R(v)$  (in particular,  $q_1 \neq p_{n-1}$ ) by remark 61 and so  $(p_i)_{0 \leq i \leq n} \cdot (q_j)_{1 \leq j \leq m}$  would be a cycle in  $G_R^s$ , that is impossible because of acyclicity of  $G_R^s$ .  $\square$

**Definition 66** (Empire). *Let  $R \in \mathbf{PPS}_{\text{DiLL}_0}$  and let  $p \in \mathcal{P}(R)$ .*

*For every switching  $s$  of  $R$ , let  $G_R^s$  be  $s$ -correction graph of  $R$ : the  $s$ -correction graph of  $R$  rooted in  $p$  is either the snipping of  $\{p\}$  in  $G_R^s$  if  $p \in \mathcal{P}^{\text{aux}}(G_R^s)$ , or  $G_R^s$  otherwise.*

*The empire of  $p$  in  $R$ , denoted by  $\epsilon_R(p)$ , is the set of  $q \in \mathcal{P}(R)$  such that  $p$  and  $q$  are connected in all the  $s$ -correction graph of  $R$  rooted in  $p$ , for every switching  $s$  of  $R$ .*

*The boundary of  $\epsilon_R(p)$  is*

$$\partial\epsilon_R(p) = \{q \in \mathcal{P}(R) \mid \exists q' \in \mathcal{P}(R) \setminus \epsilon_R(p) \text{ and } (q, q') \text{ is a path in } R \text{ of length } 1\}.$$

## 1.8 Taylor expansion

**Definition 67** (Join of two DiLL<sub>0</sub>-structures). *Let  $R, S \in \mathbf{PPS}_{\text{DiLL}_0}$  be disjoint, let  $n \in \mathbb{N}$ , let  $p_1, \dots, p_n \in \mathcal{P}^{\text{free}}(S)$  be pairwise distinct and let*



$l_1, \dots, l_n \in \mathcal{C}^{!,?}(R)$ .<sup>16</sup> The join of  $S$  in  $R$  through  $(p_1, l_1), \dots, (p_n, l_n)$  is  $R' = (\mathcal{C}', \mathcal{I}', \mathcal{D}', \mathcal{W}', \text{auxd}', \text{bc}')$  where:

•

We say then that “ $R'$  is obtained by joining  $S$  in  $R$  through  $(p_1, l_1), \dots, (p_n, l_n)$ ”.

**Definition 68.** Let  $R \in \mathbf{PS}$ . The Taylor expansion of  $R$ , denoted by  $R^*$ , is a set of  $\text{DiLL}_0$ -proof-structures defined by induction on  $\text{depth}(R) \in \mathbb{N}$  as follows:

- if  $\text{depth}(R) = 0$ , then  $R^* = \{R\}$ ;
- if  $\text{depth}(R) > 0$ , then let  $v_1, \dots, v_n$  (for some  $n \in \mathbb{N}^*$ ) be the promotion cells of  $R$  at depth 0.

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<sup>16</sup>Possibly,  $l_i = l_j$  for some  $i, j \in \{1, \dots, n\}$  with  $i \neq j$ .



## Chapter 2

# Relational semantics

Let us consider the category **Rel** of sets and relations: the Kleisli category of the comonad associated with the finite multisets functor on **Rel** is a Cartesian closed category, i.e. a denotational model for  $\lambda$ -calculus. Such an interpretation of  $\lambda$ -terms is the same as the interpretation of the Linear Logic proof-net translating the  $\lambda$ -term in the multiset based relational model of Linear Logic. This holds for both the typed and untyped case.

In  $\lambda$ -calculus, the shift from typed to untyped semantics essentially relies on the choice of a suitable object  $D$  which is reflexive, that is such that  $D \Rightarrow D$  (the exponentiation of  $D$ ) is a retract of  $D$  (i.e. there exist two morphisms  $\mathbf{abs} : (D \Rightarrow D) \rightarrow D$  and  $\mathbf{app} : D \rightarrow (D \Rightarrow D)$  such that  $\mathbf{app} \circ \mathbf{abs} = id_{D \Rightarrow D}$ ). In the multiplicative and exponential framework of Linear Logic we have more constructions than the intuitionistic arrow, then it is not enough for the object  $D$  we look for to enjoy the  $\lambda$ -calculus notion of reflexivity (it must satisfy more properties). Indeed we define an object  $D$  (definition 69) in the category **Rel** in such a way that not only  $D \times D$  and  $\mathcal{M}_{\text{fin}}(D)$  are retracts of  $D$ , but also that each of these constructs interacts well with the others (via some morphisms), thus allowing an interpretation of untyped proof-structures invariant under cut-elimination.

### 2.1 Relational spaces

We introduce a domain  $D$  to interpret (untyped) proof-structures as it is already defined in [dCPT11, dCT12]. All the following definitions are exactly the same as those ones in [dCT12].

In the definition of the domain  $D$  the set  $\{+, -\}$  of polarities is used in order to “semantically distinguish” cells of dual types  $1/\perp$ ,  $\otimes/\wp$  and  $!/?$ , which is mandatory in an untyped framework.

**Definition 69** (Atom, point). *We fix a set  $A$  not containing any pair nor any 3-tuple and such that  $*$   $\notin A$ ; we call atoms the elements of  $A$ .*

*We define  $D_n$  by induction on  $n \in \mathbb{N}$ :*

- $D_0 = A \cup (\{+, -\} \times \{*\})$ ,
- $D_{n+1} = D_0 \cup (\{+, -\} \times D_n \times D_n) \cup (\{+, -\} \times \mathcal{M}_{\text{fin}}(D_n))$ .

We set  $D = \bigcup_{n \in \mathbb{N}} D_n$ . The depth of an element  $\alpha \in D$  is the least  $n \in \mathbb{N}$  such that  $\alpha \in D_n$ .

We set  $D^{<\omega} = \bigcup_{n \in \mathbb{N}} D^n$ , whose elements are called points.

**Remark 70.**

1.  $D_n \subseteq D_{n+1}$  for every  $n \in \mathbb{N}$ . The proof is by a straightforward induction on  $n \in \mathbb{N}$ .
2. Let  $\alpha, \beta, \alpha_1, \dots, \alpha_k \in D$  (for some  $k \in \mathbb{N}$ ), let  $\gamma \in A$  and  $\iota \in \{+, -\}$ :
  - $\text{depth}(\gamma) = 0 = \text{depth}(\iota, *)$ , as  $D_0 = A \cup (\{+, -\} \times \{*\})$ ;
  - $\text{depth}(\iota, \alpha, \beta) = \max\{\text{depth}(\alpha), \text{depth}(\beta)\} + 1$ , indeed if  $\text{depth}(\alpha) = n$ ,  $\text{depth}(\beta) = m$  and  $d = \max\{n, m\}$  then  $(\alpha, \beta) \in D_d \times D_d$  and  $(\alpha, \beta) \notin D_i \times D_i$  for any  $0 \leq i \leq d-1$ , hence  $\text{depth}(\iota, \alpha, \beta) = d+1$ ;
  - $\text{depth}(\iota, [\alpha_1, \dots, \alpha_k]) = \sup\{\text{depth}(\alpha_i) \mid i \in \{1, \dots, k\}\} + 1$ , indeed if  $\text{depth}(\alpha_i) = n_i$  for any  $i \in \{1, \dots, k\}$  and  $d = \sup\{n_i \mid i \in \{1, \dots, k\}\}$ , then  $(\alpha_1, \dots, \alpha_k) \in D_d^k$  and  $(\alpha_1, \dots, \alpha_k) \notin D_j^k$  for any  $0 \leq j \leq d-1$ , hence  $\text{depth}(\iota, [\alpha_1, \dots, \alpha_k]) = d+1$ .
3. The conditions on  $A$  ensure that  $D$  satisfies the following equation

$$D = A \uplus (\{+, -\} \times \{*\}) \uplus (\{+, -\} \times D \times D) \uplus (\{+, -\} \times \mathcal{M}_{\text{fin}}(D))$$

which means that  $A$ ,  $\{+, -\} \times \{*\}$ ,  $\{+, -\} \times D \times D$  and  $\{+, -\} \times \mathcal{M}_{\text{fin}}(D)$  are retracts of  $D$ .

Thanks to remark 70.2, we can easily define some notions and prove some propositions by induction on the depth of elements of  $D$ .

The function  $()^\perp$  (which is the semantic version of the linear negation) flips polarities.

**Definition 71** (Dual). We set  $+^\perp = -$  and  $-^\perp = +$ . We define  $\alpha^\perp$  for every  $\alpha \in D$ , by induction on  $\text{depth}(\alpha) \in \mathbb{N}$  as follows (where  $\gamma \in A$ ,  $\alpha, \beta, \alpha_1, \dots, \alpha_n \in D$  for some  $n \in \mathbb{N}$ , and  $\iota \in \{+, -\}$ ):

- $\gamma^\perp = \gamma$  and  $(\iota, *)^\perp = (\iota^\perp, *)$ ;
- $(\iota, \alpha, \beta)^\perp = (\iota^\perp, \alpha^\perp, \beta^\perp)$  and  $(\iota, [\alpha_1, \dots, \alpha_n])^\perp = (\iota^\perp, [\alpha_1^\perp, \dots, \alpha_n^\perp])$ .

**Definition 72** (Substitution). A substitution is a function  $\sigma : D \rightarrow D$  induced by a function  $\sigma_A : A \rightarrow D$  and defined by induction on the depth of elements of  $D$ , as follows (where  $\gamma \in A$ ,  $\alpha, \beta, \alpha_1, \dots, \alpha_n \in D$  for some  $n \in \mathbb{N}$ , and  $\iota \in \{+, -\}$ ):

- $\sigma(\gamma) = \sigma_A(\gamma)$  and  $\sigma(\iota, *) = (\iota, *)$ ;
- $\sigma(\iota, \alpha, \beta) = (\iota, \sigma(\alpha), \sigma(\beta))$ ;
- $\sigma(\iota, [\alpha_1, \dots, \alpha_n]) = (\iota, [\sigma(\alpha_1), \dots, \sigma(\alpha_n)])$ .

If  $\sigma_A : A \rightarrow D$  is a function such that  $\text{im}(\sigma_A) \subseteq A$  (resp.  $\sigma_A$  is a bijection), then the substitution  $\sigma$  induced by  $\sigma_A$  is atomic (resp. bijective).

We denote by  $\mathfrak{M}$  (resp.  $\mathfrak{S}$ ) the set of atomic (resp. bijective and atomic) substitutions.

**Remark 73.** By a straightforward induction on  $\text{depth}(\alpha) \in \mathbb{N}$ , we can prove that  $\sigma(\alpha)^\perp = \sigma(\alpha^\perp)$  for every substitution  $\sigma$  and  $\alpha \in D$ .

**Definition 74** (Occurrences of an element of  $D$ ). For every  $\alpha \in D$ , we define  $\text{sub}(\alpha) \in \mathcal{M}_{\text{fin}}(D)$  by induction on  $\text{depth}(\alpha) \in \mathbb{N}$  as follows:

- $\text{sub}(\gamma) = [\gamma]$  if  $\gamma \in A \cup (\{+, -\} \times \{*\})$ ;
- $\text{sub}(\iota, \alpha, \beta) = [(\iota, \alpha, \beta)] + \text{sub}(\alpha) + \text{sub}(\beta)$ ;
- $\text{sub}(\iota, [\alpha_1, \dots, \alpha_n]) = [(\iota, [\alpha_1, \dots, \alpha_n])] + \sum_{j=1}^n \text{sub}(\alpha_j)$ .

For every  $n \in \mathbb{N}$  and  $(\alpha_1, \dots, \alpha_n) \in D^{<\omega}$ , we set  $\text{sub}(\alpha_1, \dots, \alpha_n) = \sum_{i=1}^n \text{sub}(\alpha_i)$ .

For every  $\alpha \in D$  and  $r \in D^{<\omega}$ , we say that  $\alpha$  occurs in  $r$  if  $\alpha \in \text{supp}(\text{sub}(r))$ , and that there are exactly  $m$  occurrences of  $\alpha$  in  $r$  if  $\text{sub}(r)(\alpha) = m$ .

In the sequel we need the notion of injective  $k$ -point of  $D^{<\omega}$  for any  $k \in \mathbb{N}$ , and for every  $E \subseteq D^{<\omega}$  the notion of  $E$ -atomic point.

**Definition 75** (Injective point,  $k$ -point,  $E$ -atomic point).  $r \in D^{<\omega}$  is injective if for every  $\gamma \in A$ , either  $\gamma$  does not occur in  $r$  or there are exactly 2 occurrences of  $\gamma$  in  $r$ . For every  $E \subseteq D^{<\omega}$ , we set  $E_{\text{inj}} = \{r \in E \mid r \text{ is injective}\}$ .

Given  $k \in \mathbb{N}$ , we say that  $r \in D^{<\omega}$  is a  $k$ -point if, for every  $m \in \mathbb{N}$  and  $\alpha_1, \dots, \alpha_m \in D$  such that  $(+, [\alpha_1, \dots, \alpha_m])$  occurs in  $r$ , we have  $m = k$ .

Let  $E \subseteq D^{<\omega}$ .  $r \in E$  is  $E$ -atomic if for every  $r' \in E$  and every substitution  $\sigma$  such that  $\sigma(r') = r$  one has  $\sigma(\gamma) \in A$  for every  $\gamma \in A$  occurring in  $r'$ . We set  $E_{\text{at}} = \{r \in E \mid r \text{ is } E\text{-atomic}\}$ .

Once the subset  $E$  of  $D^{<\omega}$  is fixed, it makes sense for  $r \in E$  to say that it is  $E$ -atomic: this means that no other element of  $E$  is “more atomic” than  $r$ . In a typed framework, we would not have to define the notion of  $E$ -atomic point, but in our untyped framework we need that: the reason will be explained after definition 80.

## 2.2 Experiments

Like in [Tor03, dCPT11, dCT12], we use experiments, introduced by Girard in [Gir87] to compute the interpretation of a proof-net in the coherent and relational semantics and deeply studied by Tortora de Falco in [Tor00, Tor03]. An experiment (definition 76) can be thought as objects between syntax and semantics allowing to associate with every ps  $R$  a point of  $D^{<\omega}$  (called result of the experiment, see definition 77) which is an element of the interpretation of  $R$  in the relational semantics. The interpretation of  $R$  in the relational semantics is the set of results of all the experiments of  $R$  (definition 78). Experiments are deeply related to non-idempotent intersection types and their derivations in the  $\lambda$ -calculus (see [dC07, dC09, Ehr12]): an experiment corresponds to a type derivation and the result of an experiment corresponds to a type. The intersection types system considered in [dC07, dC09, Ehr12] lacks idempotency and this corresponds to the fact that we use multisets for interpreting exponentials and not sets as in the set based coherent semantics introduced by Girard in [Gir87].

The definition of experiment of a ps (the same as that one in [dCT12]) is inductive and it uses the nesting condition.

**Definition 76** (Experiment). *Let  $R \in \mathbf{PS}$ . An experiment  $e$  of  $R$ , denoted by  $e : R$ , is a function associating with every  $p \in \mathcal{P}(R)$  a  $x \in \mathcal{M}_{\text{fin}}(D)$  and with every  $v \in \mathcal{C}^{\text{box}}(R)$  a finite multiset of finite multisets of experiments of  $\text{box}_R(v)$ . The definition is by induction on  $\text{depth}(R) \in \mathbb{N}$ , and we ask that  $\text{card}(e(v)) = 1$  for every  $v \in \mathcal{C}^{\text{box}}(R)$  such that  $\text{depth}_R(v) = 0$ , and  $\text{card}(e(p)) = 1$  for every  $p \in \mathcal{P}(R)$  such that  $\text{depth}_R(p) = 0$ . Furthermore the following conditions are to be fulfilled.*

1. For every  $\{p, q\} \in \mathcal{W}(R)$  such that  $\text{depth}_R(p) = 0 = \text{depth}_R(q)$ :
  - if  $\{p, q\} \in \mathcal{Ax}(R) \cup \mathcal{Cuts}(R)$ ,  $e(p) = [\alpha]$  and  $e(q) = [\beta]$ , then  $\alpha = \beta^\perp$ ;
  - if  $\{p, q\} \in \mathcal{W}(R) \setminus (\mathcal{Ax}(R) \cup \mathcal{Cuts}(R))$ , then  $e(p) = e(q)$ .
2. For every  $l \in \mathcal{C}(R)$  such that  $\text{depth}_R(\mathbf{P}_R^{\text{pri}}(l)) = 0$ :
  - if  $l \in \mathcal{C}^\otimes(R)$  (resp.  $l \in \mathcal{C}^\otimes(R)$ ),  $e(\mathbf{P}_R^{\text{left}}(l)) = [\alpha]$  and  $e(\mathbf{P}_R^{\text{right}}(l)) = [\beta]$ , then  $e(\mathbf{P}_R^{\text{pri}}(l)) = [(+, \alpha, \beta)]$  (resp.  $e(\mathbf{P}_R^{\text{pri}}(l)) = [(-, \alpha, \beta)]$ );
  - if  $l \in \mathcal{C}^1(R)$  (resp.  $l \in \mathcal{C}^\perp(R)$ ), then  $e(\mathbf{P}_R^{\text{pri}}(l)) = [(+, *)]$ , (resp.  $e(\mathbf{P}_R^{\text{pri}}(l)) = [(-, *)]$ );
  - if  $l \in \mathcal{C}^?(R)$ , then  $e(\mathbf{P}_R^{\text{pri}}(l)) = [(-, \sum_{p \in \mathcal{P}_R^{\text{aux}}(R)} e(p))]$ .
3. For every  $v \in \mathcal{C}^{\text{box}}(R)$  such that  $\text{depth}_R(\text{box}_R(v)) = 0$ , let  $e(v) = [[e_1, \dots, e_{n_v}]]$ :

- if  $p$  is the conclusion<sup>1</sup> of  $\text{box}_R(v)$  such that  $\mathbf{P}_R^{\text{pri}}(v) <_R p$ , then  $e(\mathbf{P}_R^{\text{pri}}(l)) = [(+, \sum_{i=1}^{n_v} e_i(p))]$ ;
- if  $w \in \mathcal{C}^{\text{box}}(\text{box}_R(v))$ , then  $e(w) = \sum_{i=1}^{n_v} e_i(v)$ ;<sup>2</sup>
- if  $p \in \text{inbox}_R(v)$ , then  $e(p) = \sum_{i=1}^{n_v} e_i(p)$ .<sup>3</sup>

Let  $(R, \text{ind}) \in \mathbf{PS}^{\text{ind}}$ . An experiment of  $(R, \text{ind})$  is an experiment of  $R$ .

Given a ps  $R$ , experiments of  $R$  are functions defined on  $R$  allowing to compute the interpretation of  $R$  pointwise. Indeed, for every experiment  $e$  of  $R$ , the labels associated by  $e$  with the conclusions of  $R$  form a tuple called the result of  $e$  representing which is a point of  $D^{<\omega}$ , so the result of an experiment is a truly semantic object. The set of results of all the experiments of  $R$  is the interpretation of  $R$  in the (multiset based) relational semantics.

**Definition 77** (Result of an experiment). *Let  $(R, \text{ind}) \in \mathbf{PS}^{\text{ind}}$  with  $n = \text{card}(\mathcal{P}^{\text{free}}(R))$  and let  $e$  be an experiment of  $R$ . The result of  $e$  in  $(R, \text{ind})$  is  $|e|_{\text{ind}} = (\alpha_1, \dots, \alpha_n) \in D^n$  such that  $\alpha_i$  is the unique element of the multiset  $e(\text{ind}^{-1}(i))$ , for every  $i \in \{1, \dots, n\}$ .*

**Definition 78** (Interpretation of a proof-structure). *Let  $(R, \text{ind}) \in \mathbf{PS}^{\text{ind}}$  and let  $n = \text{card}(\mathcal{P}^{\text{free}}(R))$ . The interpretation of  $(R, \text{ind})$  is*

$$\llbracket (R, \text{ind}) \rrbracket = \{|e|_{\text{ind}} \in D^n \mid e \text{ is some experiment of } R\}.$$

Experiments are defined for whatever ps, including DiLL<sub>0</sub>-ps. Any denotational semantics of DiLL<sub>0</sub>-ps provides a semantics for MELL-ps. through the Taylor expansion. This is what the following proposition says in the case of relational semantics.

**Proposition 79.** *For every  $(R, \text{ind}) \in \mathbf{PS}^{\text{ind}}$ , one has  $\llbracket (R, \text{ind}) \rrbracket = \bigcup_{\rho \in R^*} \llbracket (\rho, \text{ind}) \rrbracket$ .*

PROOF. By straightforward induction on  $\text{card}(\mathcal{C}(R)) \in \mathbb{N}$ .  $\square$

**Definition 80** (Atomic experiment). *Let  $R \in \mathbf{PS}$ . An experiment  $e$  of  $R$  is atomic if for every  $p \in \bigcup \mathcal{A}x(R)$ , one has  $e(p) \in \mathcal{M}_{\text{fin}}(A)$ .*

In our untyped framework we need to restrict the set  $E$  of all results of all experiments of a ps to the set of the results of the atomic experiments of this ps, in order to avoid the problem of “infinite  $\eta$ -expansions” which are semantically “invisible”. Of course, a given point of  $D^{<\omega}$  can be the result of an atomic experiment of a ps and the result of a non-atomic experiment of another ps. However, given  $(R, \text{ind}) \in \mathbf{PS}^{\text{ind}}$ , it makes sense for  $r \in \llbracket (R, \text{ind}) \rrbracket$  to say that it is  $\llbracket (R, \text{ind}) \rrbracket$ -atomic: this means that no other element of  $\llbracket (R, \text{ind}) \rrbracket$  is “more atomic” than  $r$ .

<sup>1</sup>See remark 42.1.

<sup>2</sup>This is well defined thanks to the nesting condition and because each promotion cell in  $\text{box}_R(v)$  is a cell of  $R$  by proposition 41, since  $\text{tr}_R(l) = ?$  for every  $l \in \mathcal{L}_0$ .

<sup>3</sup>This is well defined thanks to the nesting condition.

**Lemma 81.** *For every  $(R, \text{ind}) \in \mathbf{PS}^{\text{ind}}$  cut-free, one has*

$$\llbracket (R, \text{ind}) \rrbracket_{\text{at}} = \{|e|_{\text{ind}} \mid e \text{ is an atomic experiment of } R\}.$$

PROOF. By straightforward induction on  $\text{card}(\mathcal{C}(R)) \in \mathbb{N}$ .  $\square$

**Lemma 82.** *Let  $(R, \text{ind}) \in \mathbf{PS}^{\text{ind}}$  be cut-free and deadlock-free. For every  $r \in \llbracket (R, \text{ind}) \rrbracket_{\text{inj,at}}$ , if  $r$  is a 1-point then  $(\bar{r}, \text{ind}_r) \simeq (\text{lin}(R), \text{ind})$ .*

WHY?

PROOF. By straightforward induction on  $\text{card}(\mathcal{C}(R)) \in \mathbb{N}$ .

As  $r \in \llbracket (R, \text{ind}) \rrbracket$  is a 1-point and  $R$  is cut-free, necessarily  $\mathbf{a}_R(l) = 1$  for every  $l \in \mathcal{C}^!(R)$ .

If there exists  $l \in \mathcal{C}^{\text{term}}(R) \cap \mathcal{C}^?(R)$  and then  $\square$

Roughly speaking, lemma 82 says that, given  $(R, \text{ind}) \in \mathbf{PS}^{\text{ind}}$ , an injective and atomic 1-point in the interpretation of  $(R, \text{ind})$  is the same as  $(R, \text{ind})$  but the arrows of  $R$ .

## 2.3 The relationship between Taylor expansion and relational semantics

In the intuition of many specialists, (a result of) an experiment of a MELL-ps  $R$  is seen as a DiLL<sub>0</sub>-ps in the Taylor expansion  $R^*$  of  $R$ , and the interpretation of  $R$  in the relational semantics is seen as  $R^*$ . But this relationship between Taylor expansion and relational semantics has been never formulated precisely. This is what we aim at doing here. Quite surprisingly, the relationship between a result of an experiment of a MELL-ps  $R$  and a DiLL<sub>0</sub>-ps in  $R^*$  can be stated in the expected intuitive way only when  $R$  is cut-free. This is due to the fact seen in section ?? that two distinct DiLL<sub>0</sub>-ps in the Taylor expansion  $R^*$  of a MELL-ps  $R$  (with cuts) having the same normal form.

**Definition 83** (From points to pseudo-structures). *Let  $\alpha \in D$ . We define by induction on  $\text{depth}(\alpha)$  a pair  $(\tilde{\alpha}, ax_\alpha)$  such that  $\tilde{\alpha}$  is a pseudo-structure having only one conclusion (denoted by  $c(\alpha)$ ) and  $ax_\alpha$  is a function associating with every  $l \in \mathcal{C}^{ax}(\tilde{\alpha})$  some  $\gamma \in A$ , called the label of  $l$  as follows:*

- if  $\alpha \in A$  then  $\tilde{\alpha} = (\mathbb{C}, \emptyset, \emptyset, \emptyset, \emptyset, \emptyset)$  where  $\mathbb{C}$  is the pseudo-cell-base consisting only of a  $ax$ -cell  $l$ , and  $ax_\alpha(l) = \alpha$ ;
- if  $\alpha = (+, *)$  (resp.  $\alpha = (-, *)$ ) then  $\tilde{\alpha} = (\mathbb{C}, \emptyset, \emptyset, \emptyset, \emptyset, \emptyset)$  where  $\mathbb{C}$  is the pseudo-cell-base consisting only of a 1-cell (resp.  $\perp$ -cell), and  $ax_\alpha = \emptyset$  (the empty function);
- if  $\alpha = (+, \alpha_1, \alpha_2)$  (resp.  $\alpha = (-, \alpha_1, \alpha_2)$ ) then  $\tilde{\alpha} = (\mathbb{C}, \emptyset, \emptyset, \mathcal{W}(\tilde{\alpha}_1) \uplus \mathcal{W}(\tilde{\alpha}_2) \uplus \{w_1, w_2\}, \emptyset, \emptyset)$ , where  $\mathbb{C}$  is the pseudo-cell-base consisting of the disjoint union of  $\mathbb{C}(\tilde{\alpha}_1)$ ,  $\mathbb{C}(\tilde{\alpha}_2)$  and a  $\otimes$ -cell (resp.  $\wp$ -cell)  $l$ , and  $w_1 = \{c(\alpha_1), \mathbf{P}_{\mathbb{C}}^{\text{left}}(l)\}$  and  $w_2 = \{c(\alpha_2), \mathbf{P}_{\mathbb{C}}^{\text{right}}(l)\}$ ; moreover  $ax_\alpha = ax_{\alpha_1} \uplus ax_{\alpha_2}$ ;



- if  $\alpha = (+, [\alpha_1, \dots, \alpha_n])$  (resp.  $\alpha = (-, [\alpha_1, \dots, \alpha_n])$ ) with  $n \in \mathbb{N}$ , then  $\tilde{\alpha} = (\mathbb{C}, \emptyset, \emptyset, \bigsqcup_{0 \leq i \leq n} \mathcal{W}(\tilde{\alpha}_i) \uplus \{w_1, \dots, w_n\}, \emptyset, \emptyset)$ , where  $\mathbb{C}$  is the pseudo-cell-base consisting of the disjoint union of  $\mathbb{C}(\tilde{\alpha}_1), \dots, \mathbb{C}(\tilde{\alpha}_n)$  and a  $n$ -ary  $!$ -cell (resp.  $?$ -cell)  $l$  with  $\mathcal{P}_l^{\text{aux}}(\mathbb{C}) = \{p_1, \dots, p_n\}$ , and  $w_i = \{c(\alpha_i), p_i\}$  for every  $1 \leq i \leq n$ ; moreover  $ax_\alpha = \bigsqcup_{i=1}^n ax_{\alpha_i}$ .

For every  $n \in \mathbb{N}$  and  $r = (\alpha_1, \dots, \alpha_n) \in D^n$ , we set  $\tilde{r} = \bigsqcup_{i=1}^n \tilde{\alpha}_i$  and  $ax_r = \bigsqcup_{i=1}^n ax_{\alpha_i}$ , and we define the function  $\text{ind}_r : \{c(\alpha_1), \dots, c(\alpha_n)\} \rightarrow \{1, \dots, n\}$  by  $\text{ind}_r(c(\alpha_i)) = i$  for every  $1 \leq i \leq n$ .

Let  $r \in D^{<\omega}$ . We denote by  $\text{pseudo}(r)$  the set defined as follows:  $\Phi \in \text{pseudo}(r)$  iff  $\Phi$  is obtained from  $\tilde{r}$  by connecting  $(l_1, l'_1), \dots, (l_n, l'_n)$  (for some  $n \in \mathbb{N}$ ) where  $l_1, l'_1, \dots, l_n, l'_n$  are pairwise distinct  $ax$ -cells of  $\tilde{r}$  such that  $ax_r(l_i) = ax_r(l'_i)$ .

**Remark 84.** Let  $r \in D^{<\omega}$ .

$\tilde{r}$  (resp.  $(\tilde{r}, \text{ind}_r)$ ) is a cut-free and deadlock-free pseudo-structure (resp. indexed pseudo-structure) such that  $\mathcal{I}(\Phi) = \emptyset$  and  $\mathcal{C}^{\text{prom}}(\Phi) = \emptyset$ .

If  $r$  is injective then there exists exactly one  $\Phi \in \text{pseudo}(r)$  such that  $\Phi \in \mathbf{PPS}_{\text{DiLL}_0}$ : this is the  $\Phi \in \text{pseudo}(r)$  obtained from  $\tilde{r}$  by connecting all the pairs of distinct  $ax$ -cells with the same label in  $\tilde{r}$  (there is exactly one way to do that because of the injectivity of  $r$ ). We denote such a  $\Phi \in \text{pseudo}(r)$  by  $\bar{r}$ . Clearly,  $(\bar{r}, \text{ind}_r) \in \mathbf{PPS}_{\text{DiLL}_0}^{\text{ind}}$ .

We will see that, for every  $r \in D^{<\omega}$ ,  $\bar{r}$  can be seen as a sort of canonical representative of  $r$ .

**Proposition 85.** Let  $r, r' \in D^{<\omega}$  be injective points. If  $r \sim r'$  then  $\bar{r} \simeq \bar{r}'$ .

**Lemma 86.** Let  $(R, \text{ind}) \in \mathbf{PS}_{\text{MELL}}^{\text{ind}}$  be cut-free and deadlock-free. For every  $\rho \in (R, \text{ind})^*$  and for every atomic experiment  $e$  of  $\rho$ , if  $r$  is injective then  $(\bar{r}, \text{ind}_r) \simeq \rho$ .

PROOF. By induction on  $\text{card}(\mathcal{C}(R))$ . □

**Lemma 87.** Let  $(R, \text{ind}) \in \mathbf{PS}_{\text{MELL}}^{\text{ind}}$ . For every  $r \in \llbracket (R, \text{ind}) \rrbracket_{\text{at}, \text{inj}}$ , one has  $(\bar{r}, \text{ind}_r) \in (R, \text{ind})^*$ .

PROOF. By induction on  $\text{card}(\mathcal{C}(R))$ . □

**Theorem 88.** For every  $(R, \text{ind}) \in \mathbf{PS}_{\text{MELL}}^{\text{ind}}$ , If  $R$  is cut-free and deadlock-free then

$$(R, \text{ind})^* = \{(\bar{r}, \text{ind}_r) \mid r \in \llbracket (R, \text{ind}) \rrbracket_{\text{at}, \text{inj}}\}.$$

PROOF. Immediate consequence of lemmas 86 and 87. □

## 2.4 The connected case

**Definition 89** (Separable  $\text{?}$ -cell). *Let  $R \in \mathbf{PS}$  and let  $l \in \mathcal{C}^{\text{term}}(R) \cap \mathcal{C}^{\text{?}}(R)$ .*

*$l$  is inseparable if either  $\mathbf{a}_R(l) \leq 1$  or there exists  $v \in \mathcal{C}^{\text{box}}(R)$  such that  $\text{depth}_R(v) = 0$  and  $\mathcal{P}_l^{\text{aux}}(R) \subseteq \text{auxd}_R(v)$ .*

*$l$  is separable if it is not inseparable.*

In other words, given a ps  $R$ , a terminal  $\text{?}$ -cell in  $R$  of arity greater than or equal to 2 is inseparable if all its auxiliary ports are auxiliary doors of one and only one promotion cell at depth 0 in  $R$ .

**Remark 90.** Let  $R \in \mathbf{PS}$  and let  $l \in \mathcal{C}^{\text{term}}(R) \cap \mathcal{C}^{\text{?}}(R)$ . According to the nesting condition, if  $l$  is separable then  $\mathbf{a}_R(l) \geq 2$  and:

- either there exists  $p \in \mathcal{P}_l^{\text{aux}}(R)$  such that  $p \notin \text{Auxdoors}(R)$ ;
- or for every  $p \in \mathcal{P}_l^{\text{aux}}(R)$  there exists  $v \in \mathcal{C}^{\text{box}}(R)$  such that  $\text{depth}_R(v) = 0$  and  $p \in \text{auxd}_R(v)$ ; moreover, there exist  $p, p' \in \mathcal{P}_l^{\text{aux}}(R)$  and  $v, v' \in \mathcal{C}^{\text{box}}(R)$  such that  $v \neq v'$ ,  $\text{depth}_R(v) = 0 = \text{depth}_R(v')$  and  $p \in \text{auxd}_R(v)$  and  $p' \in \text{auxd}_R(v')$ .

We give a notion of measure of a ps.

**Definition 91** (*aux-measure*). *Let  $R \in \mathbf{PS}$  and let  $m_R^?$  be a multiset on  $\mathbb{N}$  defined by  $m_R^?(i) = \text{card}(\{l \in \mathcal{C}^{\text{?}}(R) \mid \mathbf{a}_R(l) = i\})$  for every  $i \in \mathbb{N}$ . The aux-measure of  $R$  is  $\#(R) = (m_R^?, \text{card}(\mathcal{C}(R)), \text{card}(\mathcal{P}(R)))$ .*

**Remark 92.**

1. Given  $R \in \mathbf{PS}$ ,  $m_R^?$  is a finite multiset, as  $\mathcal{C}^{\text{?}}(R)$  is a finite set.
2. For every  $R, R' \in \mathbf{PS}$ , we can establish an order relation between  $m_R^?$  and  $m_{R'}^?$ , given by the usual multiset order:

$$m_R^? \leq m_{R'}^? \Leftrightarrow \text{for any } i \in \mathbb{N}, \text{ if } m_R^?(i) > m_{R'}^?(i) \text{ then there is } j > i \text{ such that } m_R^?(j) < m_{R'}^?(j)$$

**Definition 93** (1- and 2-DiLL<sub>0</sub>-ps of a MELL-ps). *Let  $R \in \mathbf{PS}_{\text{MELL}}$  and let  $\Phi \in R^*$ .*

*$\Phi$  is the 1-(resp. 2-)DiLL<sub>0</sub>-ps of  $R$  if for every  $l \in \mathcal{C}^1(\Phi)$ , one has  $\mathbf{a}_\Phi(l) = 1$  (resp.  $\mathbf{a}_\Phi(l) = 2$ ).*

**Definition 94** (1-projection). *Let  $R \in \mathbf{PS}_{\text{MELL}}$  and let  $\Phi$  be a 2-DiLL<sub>0</sub>-ps of  $R$ . We define a 3-tuple  $(1(\Phi), \pi_{\mathcal{P}(\Phi)}, \pi_{\mathcal{C}(\Phi)})$  as follows, by induction on  $\#(\Phi)$  with the lexicographical order on  $\mathbb{N}^3$ .*

- If  $R$  is the empty ps, then  $1(\Phi)$  is the empty ps,  $\pi_{\mathcal{P}(\Phi)}$  and  $\pi_{\mathcal{C}(\Phi)}$  are the empty functions.

•

- If  $\mathfrak{a}_\Phi(l) = 0$  for any  $l \in \mathcal{C}(\Phi)$ , then  $1(\Phi) = \Phi$ ,  $\pi_{\mathcal{P}(\Phi)} = id_{\mathcal{P}(\Phi)}$  and  $\pi_{\mathcal{C}(\Phi)} = id_{\mathcal{C}(\Phi)}$ ;
- If there exists  $l \in \mathcal{C}^{\text{term}}(\Phi) \cap \mathcal{C}^{\otimes, \exists}(\Phi)$ , let  $\mathcal{P}_l^{\text{aux}}(\Phi) = \{p_l, p_r\}$  with  $p_l = \mathbf{P}_\Phi^{\text{left}}(l)$  and  $p_r = \mathbf{P}_\Phi^{\text{right}}(l)$ , let  $\Phi'$  be the erasure of  $\mathcal{P}_l(\Phi)$  in  $\Phi$ ; by induction hypothesis, one has  $\#(\Phi') < \#(\Phi)$ , so there exists  $(1(\Phi'), \pi_{\mathcal{P}(\Phi')}, \pi_{\mathcal{C}(\Phi')})$ . We set:
  - $1(\Phi)$  is the add of  $l$  in  $1(\Phi')$  in such a way that  $\mathbf{P}_{1(\Phi)}^{\text{left}}(l) = p_l$ ,  $\mathbf{P}_{1(\Phi)}^{\text{right}}(l) = p_r$  and if  $\{q, p_l\} \in \mathcal{W}(\Phi)$  (resp.  $\{q', p_r\} \in \mathcal{W}(\Phi)$ ) then  $\{q, p_l\} \in \mathcal{W}(1(\Phi))$  (resp.  $\{q', p_r\} \in \mathcal{W}(\Phi)$ );
  - $\pi_{\mathcal{P}(\Phi)} = \pi_{\mathcal{P}(\Phi')} \cup id_{\mathcal{P}_l(\Phi)}$ ;
  - $\pi_{\mathcal{C}(\Phi)} = \pi_{\mathcal{C}(\Phi')} \cup id_{\{l\}}$ .
- If there exists  $l \in \mathcal{C}^{\text{term}}(\Phi) \cap \mathcal{C}^?(\Phi)$  such that either there exists  $p \in \mathcal{P}_l^{\text{aux}}(\Phi)$  such that  $p \notin \epsilon_\Phi(l')$  for any  $l' \in \mathcal{C}^l(\Phi)$ , or there exist

**Lemma 95.** *Let  $\Phi$  be a 2-DiLL<sub>0</sub>-ps and let  $R \in \mathbf{PS}_{\text{MELL}}$  be such that  $\Phi, 1(\Phi) \in R^*$ . For every  $R' \in \mathbf{PS}_{\text{MELL}}$ , if  $\Phi, 1(\Phi) \in R'^*$  then  $R \simeq R'$ .*

PROOF. We prove by induction on that  $R = R'$  up to isomorphisms. For the sake of simplicity, we ignore in this proof all the problems related to isomorphisms.

We show □

The following theorem says that a cut-free and deadlock-free **MELL**-proof-net  $R$  is completely characterized by its atomic 2-point in its interpretation in the relational semantics.

**Theorem 96.** *Let  $r, r' \in D^{<\omega}$  be 2-points, let  $(R, \text{ind}), (R', \text{ind}') \in \mathbf{PN}^{\text{ind}}$  be cut-free, deadlock-free and such that  $r \in \llbracket (R, \text{ind}) \rrbracket_{\text{inj, at}}$  and  $r' \in \llbracket (R', \text{ind}') \rrbracket_{\text{inj, at}}$ . If  $\bar{r} = \bar{r}'$  then  $(R, \text{ind}) \simeq (R', \text{ind}')$ .*



## Part II

# Call-by-value lambda calculus



## Chapter 3

# About a call-by-value $\lambda$ -calculus

First formulated by Alonzo Church in 1936,  $\lambda$ -calculus is a formal system in mathematical logic and theoretical computer science for expressing computation by way of variable binding and substitution. It found early successes in the area of computability theory, such as a negative answer to Hilbert's Entscheidungsproblem. As pointed out by Peter Landin's 1965 paper [Lan65], sequential procedural programming languages can be understood in terms of the  $\lambda$ -calculus, which provides the basic mechanisms for procedural abstraction and procedure (subprogram) application. The  $\lambda$ -calculus may be seen as the idealized prototype of functional programming languages, like Lisp, Haskell or the various dialects of ML. Under this view,  $\beta$ -reduction (the operation performing substitution of a bound variable for an argument) corresponds to a computational step.

Because of the importance of the notion of variable binding and substitution, there is not just one system of  $\lambda$ -calculus, and in particular there are typed and untyped variants. Historically, the most important system was the untyped  $\lambda$ -calculus, in which function application has no restrictions (so the notion of the domain of a function is not built into the system). In the Church–Turing Thesis, the untyped lambda calculus is claimed to be capable of computing all effectively calculable functions; actually untyped  $\lambda$ -calculus is equivalent to all the models of computation having the highest expressive power nowadays known, like Turing machines and recursive functions. The typed  $\lambda$ -calculus is a variety that restricts function application, so that functions can only be applied if they are capable of accepting the given input's "type" of data.

Another variant of  $\lambda$ -calculus is the "call-by-value"  $\lambda$ -calculus. The most commonly used parameter passing policy for programming languages is call-by-value (CBV). Landin in [Lan65] pioneered a CBV formal evaluation for a lambda-core of ALGOL60 (named ISWIM) via the SECD abstract machine.

Ten years later, Plotkin in [Plö75] introduced the  $\lambda_{\beta_v}$ -calculus in order to grasp the CBV paradigm in a pure lambda-calculus setting. The  $\lambda_{\beta_v}$ -calculus narrows the  $\beta$ -reduction rule by allowing the reduction of a redex  $(\lambda x t)u$ , only in case  $u$  is a value, i.e. a variable or an abstraction.

### 3.1 A call-by-value $\lambda$ -calculus

We will study now  $\Lambda_{\text{CBV}}$ , a call-by-value  $\lambda$ -calculus introduced in [Ehr12] by Ehrhard and inspired by his analysis of the relational model for Linear Logic.

#### 3.1.1 The syntax of $\Lambda_{\text{CBV}}$

Let  $\mathcal{V}$  be a countable set whose elements, denoted by  $x, y, z, \dots$ , are called *variables*.

**Definition 97.** We define the elements of the sets  $\Lambda_{\text{t}}$  (terms),  $\Lambda_{\text{v}}$  (values),  $\Lambda_{\text{CBV}}$  (expressions) by mutual induction as follows:

$$\begin{array}{lll} \Lambda_{\text{t}} & L, M, N ::= (M)N \mid (V)^\dagger & \text{terms} \\ \Lambda_{\text{v}} & U, V, W ::= x \mid \lambda x M & \text{values} \\ \Lambda_{\text{CBV}} & D, E, F ::= M \mid V & \text{expressions} \end{array}$$

Note that  $\Lambda_{\text{CBV}} = \Lambda_{\text{t}} \uplus \Lambda_{\text{v}}$ . Terms of the shape  $(V)^\dagger$  (resp.  $(M)N$ ) for some value  $V$  (resp. terms  $M$  and  $N$ ) are called *promoted values* (resp. *applications*). Terms of the shape  $(M)N$  for some terms  $M$  and  $N$  are called *applications*,  $M$  (resp.  $N$ ) is in *function* (resp. *argument*) position. Values of the shape  $\lambda x M$  for some term  $M$  are called *abstractions*.

**Notation.** We follow the Krivine's notation (see [Kri93]) for applications, where the parentheses are on the function. For instance, the term  $(M)(N)L$  according to our notation is the term  $M(NL)$  according to Barendregt's notation (see [Bar84]).

Let  $M, N_1, \dots, N_n$  be terms, with  $n \in \mathbb{N}$ : if no ambiguity arise, often we use the notation  $(M)N_1 \dots N_n$  or  $MN_1 \dots N_n$  for  $(\dots((M)N_1)\dots)N_n$ , in particular if  $n = 0$  then it stands for  $M$ .

If  $n = 0$ ,  $(N_1) \dots (N_n)M$  stands for  $M$ .

If  $V$  is a value, often we write  $V^\dagger$  instead of  $(V)^\dagger$ .

**Definition 98.** With every expression  $E$  is associated its size  $\text{size}(E) \in \mathbb{N}^*$ , defined by induction on  $E$  as follows:

- $\text{size}(x) = 1$ ;
- $\text{size}(\lambda x M) = \text{size}(M) + 1$ ;
- $\text{size}(MN) = \text{size}(M) + \text{size}(N) + 1$ .
- $\text{size}(V^\dagger) = \text{size}(V) + 1$ ;



For every expression  $E$ ,  $\text{size}(E)$  is the number of rules of definition 97 used to build  $E$ , in other words it is the sum of the nodes in the tree-like representation  $\mathcal{T}_{\text{@}}M$  of  $M$ .

Due to the presence of constructor  $()^!$  (which allows to separate terms and values into two distinct sets), the set  $\Lambda_{\text{CBV}}$  of expressions does not coincide with the set  $\Lambda$  of ordinary  $\lambda$ -terms. By the way, there is an obvious “forgetful functor”  $F$  from  $\Lambda_{\text{CBV}}$  to  $\Lambda$ , defined as follows (by induction on the expression in  $\Lambda_{\text{CBV}}$ ):

$$\begin{aligned} F(x) &= x & F(\lambda x M) &= \lambda x F(M) \\ F(V^!) &= F(V) & F(MN) &= F(M)F(N) \end{aligned}$$

Definitions of free variables,  $\alpha$ -equivalence and substitution (avoiding variable capture) are extended to expressions as expected. For instance, the *free occurrences* of a variable  $x$  in an expression  $E$  are defined, by induction on  $E$ , as follows :

- if  $E$  is the variable  $x$ , then the occurrence of  $x$  in  $E$  is free;
- if  $E = (M)N$  for some terms  $M$  and  $N$ , then the free occurrences of  $x$  in  $E$  are those of  $x$  in  $M$  and  $N$ ;
- if  $E = \lambda y M$  for some term  $M$ , the free occurrences of  $x$  in  $E$  are those of  $x$  in  $M$ , except if  $x = y$ ; in that case, no occurrence of  $x$  in  $E$  is free.
- if  $E = (V)^!$  for some value  $V$ , the free occurrences of  $x$  in  $E$  are those of  $x$  in  $V$ .

A *free variable* in an expression  $E$  is a variable which has at least one free occurrence in  $E$ ; the set of free variables in  $E$  is denoted by  $\text{fv}(E)$ . An expression which has no free variable is said *closed*. A *bound variable* in an expression  $E$  is a variable which occurs in  $E$  just after the symbol  $\lambda$ . In an expression  $\lambda x M$  the  $\lambda x$  before  $M$  *binds* the free occurrences of  $x$  in  $M$ .

We work up to  $\alpha$ -equivalence.

As another example of notion coming from ordinary  $\lambda$ -calculus trivially adapted to  $\Lambda_{\text{CBV}}$ , the operation of substitution avoiding variable capture is extended by setting

$$V^![W_1/x_1, \dots, W_n/x_n] = (V[W_1/x_1, \dots, W_n/x_n])^!$$

for any values  $V, W_1, \dots, W_n$ , variables  $x_1, \dots, x_n$  and  $n \in \mathbb{N}$ . Notice that the substitution is defined only for values replacing variables. The following lemma extends at  $\Lambda_{\text{CBV}}$  a substitution lemma of ordinary  $\lambda$ -calculus.

**Lemma 99.**

1. Let  $E$  be an expression, let  $V, W$  be values and let  $x, y$  be variables. If  $x \notin \text{fv}(W) \cup \{y\}$  then  $E[V/x][W/y] = E[W/y][V[W/y]/x]$ .

2. If  $E$  is a vale (resp. a term),  $V$  is a value and  $x$  is a variable, then  $E[V/x]$  is a value (resp. a term).

PROOF.

1. By induction on the expression  $E$ . The only novelty is the case where  $E = U^!$  for some value  $U$ : by applying induction hypothesis the identity holds.
2. By a straightforward induction on the expression  $E$ . □

**Remark 100.** It follows immediately from definition that:

1. Every term is in the form

$$(V^!)M_1 \dots M_m \quad \text{and} \quad (N_1) \dots (N_n)W^!$$

where  $V$  and  $W$  are values,  $m, n \in \mathbb{N}$ ,  $M_i$  and  $N_j$  are terms for every  $1 \leq i \leq m$  and  $1 \leq j \leq n$ ; both of these forms are unique, and  $m = 0$  iff  $n = 0$ ; moreover if  $m = 0$  (i.e.  $n = 0$ ) then  $V = W$ , otherwise  $M_m = (N_2) \dots (N_n)W^!$  and  $N_1 = (V^!)M_1 \dots M_{m-1}$ .

2. As any term has always a finite length, applying recursively the left (resp. right) decomposition in remark 100.1 yields that every term  $M$  either is such that  $M = V^!$  for some value  $V$  or there exist  $\ell \in \mathbb{N}$ , values  $V, V_0, \dots, V_\ell$  and terms  $L_0, \dots, L_\ell$  such that:

- $L_0 = (V^!)V_0^!L_{01} \dots L_{0k_0}$  (resp.  $L_0 = (L_{01}) \dots (L_{0k_0})(V^!)V_0^!$ ) for some  $k_0 \in \mathbb{N}$  and terms  $L_{01}, \dots, L_{0k_0}$ ;
- for every  $1 \leq i \leq \ell$ , we have  $L_i = (V_i^!)L_{i-1}L_{i1} \dots L_{ik_i}$  (resp.  $L_i = (L_{i1}) \dots (L_{ik_i})(L_{i-1})V_i^!$ ) for some  $k_i \in \mathbb{N}$  and terms  $L_{i1}, \dots, L_{ik_i}$ ;
- $M = L_\ell$ .

Both of this decompositions are unique. It is more natural to consider these decompositions as a binary tree, see §3.4.

3. As any term has always a finite length and any value is either a variable or of the shape  $\lambda xM$  for some term  $M$ , applying recursively the left (resp. right) decomposition in remark 100.2 yields that every term  $M$  either is such that  $M = y^!$  for some variable  $y$  or there exist  $m, \ell_0, \dots, \ell_m \in \mathbb{N}$ , terms  $M_0, \dots, M_m$ , a variable  $y$  and for every  $0 \leq i \leq m$  values  $V_{i0}, \dots, V_{i\ell_i}$  and terms  $L_{i0}, \dots, L_{i\ell_i}$ , such that:

- $L_{00} = (y^!)V_{00}^!L_{001} \dots L_{00k_{00}}$  (resp.  $L_{00} = (L_{001}) \dots (L_{00k_{00}})(y^!)V_{00}^!$ ) for some  $k_{i0} \in \mathbb{N}$  and terms  $L_{i01}, \dots, L_{i0k_0}$ ;

WHY?

- for every  $1 \leq j \leq \ell_i$ , we have  $L_{ij} = (V_{ij}^!)L_{ij-1}L_{ij1} \cdots L_{ijk_{ij}}$  (resp.  $L_{ij} = (L_{ij1}) \cdots (L_{ijk_{ij}})(L_{ij-1})V_j^!$ ) for some  $k_j \in \mathbb{N}$  and terms  $L_{j1}, \dots, L_{jk_j}$ ;
- $M_i = L_{i\ell_i}$ .
- for every  $1 \leq i \leq m$ :
  - $L_{i0} = ((\lambda x_i M_{i-1})^!)V_{i0}^!L_{i01} \cdots L_{i0k_{i0}}$  (resp.  $L_{i0} = (L_{i01}) \cdots (L_{i0k_0})(\lambda x_i M_{i-1})^!V_{i0}^!$ ) for some  $k_{i0} \in \mathbb{N}$  and terms  $L_{i01}, \dots, L_{i0k_0}$ ;
  - for every  $1 \leq j \leq \ell_i$ , we have  $L_{ij} = (V_{ij}^!)L_{ij-1}L_{ij1} \cdots L_{ijk_{ij}}$  (resp.  $L_{ij} = (L_{ij1}) \cdots (L_{ijk_{ij}})(L_{ij-1})V_j^!$ ) for some  $k_j \in \mathbb{N}$  and terms  $L_{j1}, \dots, L_{jk_j}$ ;
  - $M_i = L_{i\ell_i}$ .
- 

Both of these decompositions are unique. It is more natural to consider these decompositions as an iteration of binary trees, see §3.4.

4. Every closed value is in the form  $\lambda x N$ , where  $N$  is a term with  $\text{fv}(N) \subseteq \{x\}$ .

**Definition 101.** For every expression  $E$ , we define by induction on  $E$  the set  $\text{sub}(E)$  of the subexpressions of  $E$  as follows:

$$\begin{array}{ll} \text{sub}(x) = \{x\} & \text{sub}(\lambda x M) = \text{sub}(M) \cup \{\lambda x M\} \\ \text{sub}(V^!) = \text{sub}(V) \cup \{V^!\} & \text{sub}(MN) = \text{sub}(M) \cup \text{sub}(N) \cup \{(M)N\} \end{array}$$

A subterm (resp. subvalue) of an expression  $E$  is a term (resp. value) which is a subexpression of  $E$ .

**Definition 102** ( $\beta_v$ - and  $\hat{\beta}_v$ -redex). A  $\beta$ -redex is a term of the shape  $(\lambda x M)^!N$  for some terms  $M$  and  $N$ . A  $\beta_v$ -redex is a term of the shape  $(\lambda x M)^!V^!$  for some value  $V$  and term  $M$ , its contractum is the term  $M[V/x]$ .

A  $\sigma_1$  (resp.  $\sigma'_3$ )-redex is a term of the shape  $(\lambda x M)^!NL$  (resp.  $(M)(\lambda x L)^!N$ ) for some terms  $M$ ,  $N$  and  $L$  with  $x \notin \text{fv}(L)$  (resp.  $x \notin \text{fv}(M)$ ), its contractum is the term  $(\lambda x ML)^!N$ . A  $\sigma_3$ -redex is a  $\sigma'_3$ -redex  $(M)(\lambda x L)^!N$  such that  $M = V^!$  for some value  $V$ .

A  $\sigma_v$  (resp.  $\sigma'_v$ )-redex is either a  $\sigma_1$ -redex or a  $\sigma_3$  (resp.  $\sigma'_3$ )-redex. A  $\beta_{v\sigma}$  (resp.  $\beta_{v\sigma'}$ )-redex is either a  $\beta_v$ -redex or a  $\sigma_v$  (resp.  $\sigma'_v$ )-redex.

Let  $E$  be an expression and let  $R \in \{\beta, \beta_v, \sigma_1, \sigma_3, \sigma'_3, \sigma_v, \sigma'_v, \beta_{v\sigma}, \beta_{v\sigma'}\}$ . A  $R$ -redex in  $E$  is an occurrence in  $E$  of a subterm of  $E$  which is a  $R$ -redex. A  $\hat{R}$ -redex in  $E$  is a  $R$ -redex in  $E$  which is not in any subvalue of  $E$ . We say that  $E$  contains a  $R$  (resp.  $\hat{R}$ )-redex if there is a  $R$  (resp.  $\hat{R}$ )-redex in  $E$ .

In order to compare  $\beta_v$ -redexes in  $\Lambda_{\text{CBV}}$  with  $\beta$ -redexes in ordinary (call-by-name)  $\lambda$ -calculus by means of the “forgetful functor”  $F$ , we observe that if a

term  $M \in \Lambda_{\text{CBV}}$  is a  $\beta_v$ -redex then  $F(M)$  is a  $\beta$ -redex (in  $\Lambda$ ), but the converse does not hold: for instance,  $M = (\lambda x x^!)(y^!)z^!$  is not a  $\beta_v$ -redex, differently from  $F(M) = (\lambda x x)(y)z$  which is a  $\beta$ -redex. Indeed for every term  $t \in \Lambda$  there exists a  $\beta_v$ -redex  $M \in \Lambda_{\text{CBV}}$  such that  $F(M) = t$  iff  $t = (\lambda x u)v$  for some  $u, v \in \Lambda$  such that  $v$  is a variable or an abstraction. Essentially, modulo the “forgetful functor”  $F$ , a  $\beta_v$ -redex is a  $\beta$ -redex such that its argument is a value.

A  $\hat{\beta}_v$ -redex can be seen as a “outermost”  $\beta_v$ -redex, that is a  $\beta_v$ -redex not contained in any other  $\beta_v$ -redex.

**Remark 103.**

1. Every term of the shape  $V^!$  for some value  $V$  contains no  $\hat{\beta}_v$ -redexes, indeed a variable is not a  $\hat{\beta}_v$ -redex and, for every term  $N$ , any possible  $\hat{\beta}_v$ -redex in  $N$  is “invisible” in  $(\lambda x N)^!$  since  $\lambda x N$  is a subvalue of  $(\lambda x N)^!$ .  
On the contrary, a value  $V$  (and so a term  $V^!$ ) might contains a  $\beta_v$ -redex. For instance  $(\lambda y x^!)^!x^!$  is a  $\beta_v$ -redex in the value  $\lambda d(\lambda y x^!)^!z^!$  (and in the term  $(\lambda d(\lambda y x^!)^!z^!)^!$ ).
2. If a term contains several  $\hat{\beta}_v$ -redexes then they are non-overlapping, i.e. they have no common occurrences of subexpressions. Indeed, a  $\hat{\beta}_v$ -redex is of the shape  $(\lambda x N)^!V^!$  for some term  $N$  and value  $V$ , where  $(\lambda x N)^!$  and  $V^!$  contain no  $\hat{\beta}_v$ -redexes (see remark 103.1).

The following definitions will be used to define and characterize binary relations on the set  $\Lambda_{\text{CBV}}$  of expressions.

**Definition 104** (Contextual and applicative closure). *Let  $R$  be a binary relation on  $\Lambda_{\text{CBV}}$ .*

*We say that  $R$  passes to context (resp.  $R$  passes to applicative contexts) if  $R$  is such that the following conditions 1, 2, 3 and 4 (resp. 1 and 2) hold, for any terms  $M, M', N$  and values  $V, V'$ :*

1. *if  $M R M'$  then  $MN R M'N$ ;*
2. *if  $M R M'$  then  $NM R NM'$ ;*
3. *if  $V R V'$  then  $V^! R V'^!$ ;*
4. *if  $M R M'$  then  $\lambda x M R \lambda x M'$ .*

*The contextual closure of  $R$  is the binary relation  $R'$  on  $\Lambda_{\text{CBV}}$  defined by applying, a finite number of times, the following rules:*

$$\frac{M R N}{M R' N} R \quad \frac{M R' M'}{MN R' M'N} @_l \quad \frac{N R' N'}{MN R' MN'} @_r$$

$$\frac{M R' M'}{\lambda x M R' \lambda x M'} \lambda \quad \frac{V R' W}{V^! R' W^!} !$$

The applicative closure of  $R$  is the binary relation  $R'$  on  $\Lambda_{\text{CBV}}$  defined by applying, a finite number of times, the following rules:

$$\frac{M R N}{M R' N} R \quad \frac{M R' M'}{MN R' M'N} @_l \quad \frac{N R' N'}{MN R' MN'} @_r$$

In the sequel we will consider contextual or applicative closures  $R'$  of relations  $R$  defined only by axiom rules. Therefore, thanks to the  $R$ -rule, we are entitled to talk about the axiom rules of  $R$  as derivation rules of the relation  $R'$ .

**Notation.** Let  $R$  be a binary relation on a set  $X$ .

We denote by  $R^=$  (resp.  $R^+$ ;  $R^*$ ;  $R^T$ ) the reflexive (resp. transitive; reflexive-transitive; symmetric) closure of  $R$ . We denote by  $\simeq_R$  the symmetric and reflexive-transitive closure of  $R$ , i.e.  $\simeq_R = (R^T)^*$ .

Let  $E, F \in X$  and  $n \in \mathbb{N}$ : we say that  $E$   $R$ -reduces to  $F$  in  $n$  steps (and we write  $E R^n F$ ) if there exists a finite sequence  $(E_i)_{0 \leq i \leq n}$  of elements of  $X$  such that  $E = E_0$ ,  $F = E_n$  and  $E_i R E_{i+1}$  for every  $0 \leq i < n$ .

**Remark 105.**

1. If  $R \subseteq \Lambda_t \times \Lambda_t$  and  $R'$  is the contextual (resp. applicative) closure of  $R$ , then  $R' \subseteq (\Lambda_t \times \Lambda_t) \cup (\Lambda_v \times \Lambda_v)$  (resp.  $R' \subseteq \Lambda_t \times \Lambda_t$ ). The proof is by a straightforward induction on the derivation of  $E R' F$ , where  $E$  and  $F$  are expressions.
2. If  $R$  is a binary relation on  $\Lambda_{\text{CBV}}$  passing to context (resp. passing to applicative context) then  $R^=$ ,  $R^+$ ,  $R^*$ ,  $R^T$  and  $\simeq_R$  pass to context (resp. pass to applicative context).

We recall some standard definitions in term rewriting systems.

**Definition 106.** Let  $R$  be a binary relation on a set  $X$  and let  $E \in X$ .

$E$  is a  $R$ -normal form or is  $R$ -normal if there is no expression  $E'$  such that  $E R E'$ .

A  $R$ -normal form of  $E$  is a  $R$ -normal form  $E'$  such that  $E \beta_v^* E'$ .

$E$  is  $R$ -normalizable if there exists a  $R$ -normal form of  $E$ .

$E$  is  $R$ -strongly normalizable if there is no infinite sequence  $(E_i)_{i \in \mathbb{N}}$  of elements of  $X$  such that  $E_0 = E$  and  $E_i R E_{i+1}$  for every  $i \in \mathbb{N}$ .

**Definition 107.** Let  $R$  be a binary relation on a set  $X$ .

$R$  is strongly (resp. locally) confluent if for every  $E, E_1, E_2 \in X$  such that  $E R E_i$  for  $i \in \{1, 2\}$  there exists  $E'$  such that  $E_i R E'$  (resp.  $E_i R^* E'$ ) for  $i \in \{1, 2\}$ .

$R$  is confluent if  $R^*$  is strongly confluent.

We recall a well-known result of term rewriting system.

**Theorem 108.** Let  $R$  be a confluent binary relation on a set  $X$  and let  $E_1, E_2 \in X$ . If  $E_1 \simeq_R E_2$  then there exists  $E \in X$  such that  $E_1 R^* E$  and  $E_2 R^* E$ .

### 3.1.2 Some call-by-value $\beta$ -reductions

The following notions of  $\beta_v$ - and  $\hat{\beta}_v$ -reduction are introduced by Ehrhard in [Ehr12]. They formulate respectively the the well-known ([Plo75]) call-by-value and lazy (or weak) call-by-value  $\beta$ -reduction for the syntax presented in §3.1.1. Some results of this section are nothing but a reformulation in  $\Lambda_{\text{CBV}}$  of well-known results for call-by-value  $\lambda$ -calculus, the novelty is in pointing out the deep symmetries in  $\beta_v$ - and especially  $\hat{\beta}_v$ -reduction.

**Definition 109** ( $\beta_v$ - and  $\hat{\beta}_v$ -reduction). *The  $\beta_v$ -reduction (resp. weak  $\beta_v$ -reduction or  $\hat{\beta}_v$ -reduction), denoted by  $\beta_v$  (resp.  $\hat{\beta}_v$ ), is the contextual (resp. applicative) closure of the binary relation  $\rightarrow_{\beta_v}$  on  $\Lambda_{\text{t}}$  defined by the following rule:*

$$\frac{}{(\lambda x M)!V! \rightarrow_{\beta_v} M[V/x]} \beta$$

where  $M$  is a term and  $V$  is a value.

**Remark 110.** By remark 105,  $\beta_v \subseteq (\Lambda_{\text{t}} \times \Lambda_{\text{t}}) \cup (\Lambda_{\text{v}} \times \Lambda_{\text{v}})$  and  $\hat{\beta}_v \subseteq \Lambda_{\text{t}} \times \Lambda_{\text{t}}$ .

In order to compare our call-by-value  $\lambda$ -calculus with ordinary (call-by-name)  $\lambda$ -calculus with respect to reductions by means of the “forgetful functor”  $F$  (see p. 65), we can prove, by straightforward induction on  $E \in \Lambda_{\text{CBV}}$ , that if  $E$  and  $E'$  are expressions such that  $E \beta_v E'$  then  $F(E) \beta F(E')$ , but the converse does not hold: for instance  $M = (\lambda x x')(y!)z!$  is  $\beta_v$ -normal (in  $\Lambda_{\text{CBV}}$ ), on the contrary  $F(M) = (\lambda x x)(y)z \beta (y)z$  (in  $\Lambda$ ). In other words modulo the “forgetful functor”  $F$ , the call-by-value  $\lambda$ -calculus allows to reduce a  $\beta$ -redex only if its argument is a value, i.e. a variable or an abstraction (see the  $\beta$ -rule for  $\beta_v$ - and  $\hat{\beta}_v$ -reductions), whilst there is no such a restriction in ordinary  $\lambda$ -calculus.

**Remark 111.**

1. It is immediate to check that for every expression  $E$  (resp. term  $M$ ), there exists an expression  $E'$  (resp. a term  $M'$ ) such that  $E \beta_v E'$  (resp.  $M \hat{\beta}_v M'$ ) iff  $E$  (resp.  $M$ ) contains a  $\beta_v$  (resp.  $\hat{\beta}_v$ )-redex. Therefore, an expression (resp. a term) is  $\beta_v$  (resp.  $\hat{\beta}_v$ )-normal iff it contains no  $\beta_v$  (resp.  $\hat{\beta}_v$ )-redex.
2. It is easy to verify that for all expressions  $E, E'$  (resp. terms  $M, M'$ ),  $E \beta_v E'$  (resp.  $M \hat{\beta}_v M'$ ) iff  $E'$  (resp.  $M'$ ) is obtained from  $E$  (resp.  $M$ ) by replacing exactly one  $\beta_v$ - (resp.  $\hat{\beta}_v$ -)redex in  $E$  (resp.  $M$ ) with its contractum.
3. Clearly,  $\hat{\beta}_v \subseteq \beta_v$  (the proof is by induction on the length of the derivation of  $M \hat{\beta}_v M'$ ). More precisely, weak  $\beta_v$ -reduction is the  $\beta_v$ -reduction with the restriction that it does not reduce under the  $\lambda$ 's (whence the

word “weak”):  $\hat{\beta}_v$ -reduction reduces a  $\beta_v$ -redex only if there is no  $\lambda$  in front of it. In particular, every  $\beta_v$ -normal form is  $\hat{\beta}_v$ -normal; the converse fails to hold: for instance,  $(\lambda d(\lambda y x^!)z^!)^!$  is  $\hat{\beta}_v$ -normal but not  $\beta_v$ -normal since  $(\lambda d(\lambda y x^!)z^!)^! \beta_v (\lambda x z^!)^!$ .

4. Terms of the shape  $V^!$  where  $V$  is a value are  $\hat{\beta}_v$ -normal forms; on the contrary, a value  $V$  and so a term  $V^!$  are not necessarily  $\beta_v$ -normal (see remarks 103.1 and 111.1)
5. All the critical pairs for  $\hat{\beta}_v$ -reduction (i.e. terms  $M, M_1, M_2$  such that  $M \hat{\beta}_v M_1$  and  $M \hat{\beta}_v M_2$  with  $M_1 \neq M_2$ ) arise from non-overlapping  $\hat{\beta}_v$ -redexes in the same term (see remarks 103.2 and 111.2)
6. For every expressions  $E$  and  $E'$ , if  $E \beta_v E'$  then  $\text{fv}(E') \subseteq \text{fv}(E)$  (the proof is by a straightforward induction on  $E$ ). In particular, for every closed expression  $E$ , if  $E \beta_v E'$  then  $E'$  is closed.

In [Ehr12] Ehrhard showed that  $\beta_v$ -reduction is confluent and that  $\hat{\beta}_v$ -reduction enjoys the following propriety: a term is  $\hat{\beta}_v$ -normalizable iff its interpretation in the relational model for  $\Lambda_{\text{CBV}}$  defined in [Ehr12] is not empty. The latter result is analogous to that one in ordinary  $\lambda$ -calculus stating that a (ordinary) term is head-normalizable iff its interpretation in the Engler model is empty. This allows to draw a parallel between  $\hat{\beta}_v$ -reduction and head reduction in ordinary  $\lambda$ -calculus. The most apparent difference between these two things is that  $\hat{\beta}_v$ -reduction is not a reduction strategy, that is a term in  $\Lambda_{\text{CBV}}$  might contains several  $\hat{\beta}_v$ -redexes, whilst every term in ordinary  $\lambda$ -calculus can have at most one head redex. We will show that this is only a seeming difference.

The following notions of size will be used several times, they are well-defined for all terms by remark 100.1.

**Definition 112.** For every term  $M$ , their sizes  $\#_l M \in \mathbb{N}$  and  $\#_r M \in \mathbb{N}$  are defined by induction on  $M$  as follows:

$$\#_l M = \begin{cases} 0 & \text{if } M = V^! \text{ for some value } V; \\ 0 & \text{if } M = (V^!)W^!N_1 \dots N_n \text{ for some } n \in \mathbb{N}, \text{ terms } N_1, \dots, N_n \text{ and} \\ & \text{values } V, W; \\ \#_l L_1 + \#_l L_2 + 1 & \text{if } M = ((V^!)(L_1)L_2)N_1 \dots N_n \text{ for some } n \in \mathbb{N}, \text{ terms } N_1, \dots, N_n, \\ & L_1, L_2 \text{ and value } V. \end{cases}$$

$$\#_r M = \begin{cases} 0 & \text{if } M = V^! \text{ for some value } V; \\ 0 & \text{if } M = (N_1) \dots (N_n)(W^!)V^! \text{ for some } n \in \mathbb{N}, \text{ terms } N_1, \dots, N_n \text{ and} \\ & \text{values } V, W; \\ \#_r L_1 + \#_r L_2 + 1 & \text{if } M = (N_1) \dots (N_n)((L_1)L_2)V^! \text{ for some } n \in \mathbb{N}, \text{ terms } N_1, \dots, N_n, \\ & L_1, L_2 \text{ and value } V. \end{cases}$$

The closed  $\hat{\beta}_v$ -normal forms are promoted values easily characterizable.

**Proposition 113.** *Let  $M$  be a closed term:  $M$  is a  $\hat{\beta}_v$ -normal form iff  $M = (\lambda x N)^!$  for some term  $N$  with  $\text{fv}(N) \subseteq \{x\}$ .*

PROOF.

$\Leftarrow$ : Trivial (it is not necessary to suppose  $M$  is closed, see also remark 111.4).

$\Rightarrow$ : The proof is by induction on the size  $\#_1 M \in \mathbb{N}$ . By remark 100.1,  $M = (V^!)M_1 \dots M_m$  for some  $m \in \mathbb{N}$ , terms  $M_1, \dots, M_m$  and value  $V$ . As  $M$  is closed,  $V$  is a closed value, thus  $V = \lambda x N$  for some term  $N$  with  $\text{fv}(N) \subseteq \{x\}$ .

If  $\#_1 M = 0$  then  $m = 0$  and so  $M = (\lambda x N)^!$ , otherwise it should be  $m > 0$  and  $M_1 = W^!$  for some value  $W$  and so  $M = (\lambda x N)^! W^! M_2 \dots M_m$ , that is impossible because  $M$  is  $\hat{\beta}_v$ -normal.

If  $\#_1 M > 0$  then it should be  $m > 0$  and  $M_1 = (L_1)L_2$  for some closed terms  $L_1, L_2$  which are  $\hat{\beta}_v$ -normal forms (since  $M$  is a  $\hat{\beta}_v$ -normal form), hence  $L_1 = (\lambda x_1 N_1)^!$  and  $L_2 = (\lambda x_2 N_2)^!$  for some terms  $N_1, N_2$  by induction hypothesis, thus  $M = ((V^!)((\lambda x_1 N_1)^!)(\lambda x_2 N_2)^!)M_2 \dots M_m$ , that is impossible because  $M$  is  $\hat{\beta}_v$ -normal.

Therefore the only possibility is that  $M = (\lambda x N)^!$  for some term  $N$  with  $\text{fv}(N) \subseteq \{x\}$ . □

**Theorem 114** (Strong confluence for  $\hat{\beta}_v$ ). *Let  $M, M_1, M_2$  be terms: if  $M \hat{\beta}_v M_1$  and  $M \hat{\beta}_v M_2$  with  $M_1 \neq M_2$ , then there exists a term  $N$  such that  $M_1 \hat{\beta}_v N$  and  $M_2 \hat{\beta}_v N$ .*

PROOF. By induction on the term  $M$ . Let us consider the last rule of the derivation of  $M \hat{\beta}_v M_1$ .

If it is the  $\beta$ -rule, then  $M = (\lambda x N_1)^! V^!$  and  $M_1 = N[V/x]$ , so there is no  $M_2 \neq M_1$  such that  $M \hat{\beta}_v M_2$ , since  $(\lambda x N_1)^!$  and  $V^!$  are  $\hat{\beta}_v$ -normal forms (see remark 111.4).

If it is the  $@_l$ -rule, then  $M = N_1 N_2$  and  $M_1 = N'_1 N_2$  with  $N_1 \hat{\beta}_v N'_1$ , hence  $N_1 \neq (\lambda x M')^!$  for any term  $M'$  (see remark 111.4). Thus there are only two cases : either  $M_2 = N_1 N'_2$  with  $N_2 \hat{\beta}_v N'_2$  and then  $M_1 \hat{\beta}_v N$  and  $M_2 \hat{\beta}_v N$  where  $N = N'_1 N'_2$ ; or  $M_2 = N''_1 N_2$  with  $N_1 \hat{\beta}_v N''_1 \neq N'_1$  by hypothesis, and then there exists a term  $L$  such that  $N'_1 \hat{\beta}_v L$  and  $N''_1 \hat{\beta}_v L$  by induction hypothesis, so  $M_1 = N'_1 N_2 \hat{\beta}_v N$  and  $M_2 = N''_1 N_2 \hat{\beta}_v N$  where  $N = L N_2$ .

If it is the  $@_r$ -rule, then  $M = N_1 N_2$  and  $M_1 = N_1 N'_2$  with  $N_2 \hat{\beta}_v N'_2$ , hence  $N_2 \neq V^!$  for any value  $V$  (see remark 111.4). Thus there are only two cases: either  $M_2 = N'_1 N_2$  with  $N_1 \hat{\beta}_v N'_1$  and then  $M_1 \hat{\beta}_v N$  and  $M_2 \hat{\beta}_v N$



where  $N = N'_1 N'_2$ ; or  $M_2 = N_1 N''_2$  with  $N_2 \hat{\beta}_v N''_2 \neq N'_2$  by hypothesis, and then there exists a term  $L$  such that  $N'_2 \hat{\beta}_v L$  and  $N''_2 \hat{\beta}_v L$  by induction hypothesis, so  $M_1 = N_1 N'_2 \hat{\beta}_v N$  and  $M_2 = N_1 N''_2 \hat{\beta}_v N$  where  $N = N_1 L$ .  $\square$

The following corollary of theorem 114 is a well known result which holds for every strongly confluent term rewriting system.

**Corollary 115** (Confluence, uniqueness of normal form, number of steps).

1.  $\hat{\beta}_v$  is confluent. More precisely, let  $M, M_1, M_2$  be terms: if  $M \hat{\beta}_v^* M_1$  in  $m_1 \in \mathbb{N}$  steps and  $M \hat{\beta}_v^* M_2$  in  $m_2 \in \mathbb{N}$  steps, then there exists a term  $N$  such that  $M_1 \hat{\beta}_v^* N$  in  $n_1 \leq m_2$  steps and  $M_2 \hat{\beta}_v^* N$  in  $n_2 \leq m_1$  steps.
2. Every term  $M$  has at most a  $\hat{\beta}_v$ -normal form, and if that exists then all the  $\hat{\beta}_v$ -reductions from  $M$  to its  $\hat{\beta}_v$ -normal form have the same number of steps.
3. Every term  $M$  is  $\hat{\beta}_v$ -strongly normalizable iff it is  $\hat{\beta}_v$ -normalizable.

PROOF.

1. By induction on  $m_1 + m_2 \in \mathbb{N}$ .

If  $m_2 = 0$  then  $M = M_2$ , hence  $M_2 \hat{\beta}_v^* M_1$  in  $m_1$  steps (and  $M_1 \hat{\beta}_v^* M_1$  in 0 steps).

If  $m_1 = 0$  then  $M = M_1$ , hence  $M_1 \hat{\beta}_v^* M_2$  in  $m_2$  steps (and  $M_2 \hat{\beta}_v^* M_2$  in 0 steps).

If  $m_1, m_2 > 0$  then there exist terms  $L_1, L_2$  such that  $M \hat{\beta}_v L_1$  and  $M \hat{\beta}_v L_2$ : by theorem 114, there exist a term  $L$  such that  $L_1 \hat{\beta}_v^* L$  and  $L_2 \hat{\beta}_v^* L$  in at most one step. By induction hypothesis (as  $L_2 \hat{\beta}_v^* M_2$  in  $m_2 - 1$  steps), there exists a term  $N'$  such that  $L \hat{\beta}_v^* N'$  in  $\ell \leq m_2 - 1$  steps and  $M_2 \hat{\beta}_v^* N'$  in at most one step. Therefore  $L_1 \hat{\beta}_v^* M_1$  in  $m_1 - 1$  steps and  $L_1 \hat{\beta}_v^* N'$  in  $\ell' \leq \ell + 1 \leq m_2$  steps, so there exists a term  $N$  such that  $M_1 \hat{\beta}_v^* N$  in  $n_1 \leq \ell' \leq m_2$  steps and  $N' \hat{\beta}_v^* N$  in  $n \leq m_1 - 1$  steps by induction hypothesis, thus  $M_2 \hat{\beta}_v^* N$  in  $n_2 \leq n + 1 \leq m_1$  steps.

2. If  $M \hat{\beta}_v^* M_1$  and  $M \hat{\beta}_v^* M_2$  where  $M_1$  and  $M_2$  are  $\hat{\beta}_v$ -normal forms, then there exists a term  $N$  such that  $M_1 \hat{\beta}_v^* N$  and  $M_2 \hat{\beta}_v^* N$  by corollary 115.1, so  $M_1 = N = M_2$  since  $M_1$  and  $M_2$  are  $\hat{\beta}_v$ -normal forms.

Let  $M'$  be the  $\hat{\beta}_v$ -normal form of  $M$ . We prove by induction on  $m \in \mathbb{N}$  that if  $M \hat{\beta}_v^* M'$  in  $m$  steps, then every  $\hat{\beta}_v$ -reduction from  $M$  to  $M'$  has length  $m$ .

- If  $m = 0$  then  $M = M'$  and so  $M$  is a  $\hat{\beta}_v$ -normal form, hence the  $\hat{\beta}_v$ -reduction of 0 steps is the only  $\hat{\beta}_v$ -reduction from  $M$  to  $M'$ .

- If  $m > 0$  then there exists a term  $M_1$  such that  $M \hat{\beta}_v M_1$  and  $M_1 \hat{\beta}_v^* M'$  in  $m - 1$  steps. We show that for every term  $M_2$  and  $m_2 \in \mathbb{N}$ , if  $M \hat{\beta}_v^* M_2$  in  $m_2$  steps and  $M_2 \hat{\beta}_v^* M'$  then  $M_2 \hat{\beta}_v^* M'$  in  $m'$  steps with  $m = m_2 + m'$ .

If  $m_2 = 0$  then  $M_2 = M \hat{\beta}_v^* M'$  in  $m$  steps by hypothesis, so we conclude by taking  $m' = m$ .

If  $m_2 > 0$  then there exists a term  $N$  such that  $M \hat{\beta}_v N$  and  $N \hat{\beta}_v^* M_2$  in  $m_2 - 1$  steps. If  $N = M_1$  then  $N \hat{\beta}_v^* M'$  in  $m - 1$  steps by hypothesis, so  $M_2 \hat{\beta}_v^* M'$  in  $m'$  steps with  $m - 1 = m_2 - 1 + m'$  by induction hypothesis applied to  $M_1$ , thus  $m = m_2 + m'$ . If  $N \neq M_1$  then there exists a term  $N'$  such that  $M_1 \hat{\beta}_v N'$  and  $N \hat{\beta}_v N'$  by theorem 114, hence  $N' \hat{\beta}_v^* M'$  in  $n$  steps with  $n + 1 = m - 1$  by induction hypothesis applied to  $M_1$ , so  $N \hat{\beta}_v^* M'$  in  $n + 1 = m - 1$  steps, therefore by applying the induction hypothesis to  $N$  we conclude that  $M_2 \hat{\beta}_v^* M'$  in  $m'$  steps with  $m - 1 = m_2 - 1 + m'$ , hence  $m = m_2 + m'$ .

3. The left-to-right direction is obvious. For the right-to-left direction, let us suppose by absurd that there exists a  $\hat{\beta}_v$ -normalizable term  $M$  which is not  $\hat{\beta}_v$ -strongly normalizable: then there should exist the  $\hat{\beta}_v$ -normal form  $M'$  of  $M$  and an infinite sequence of terms  $(M_i)_{i \in \mathbb{N}}$  such that  $M = M_0$  and  $M_i \hat{\beta}_v M_{i+1}$ ; if  $M \hat{\beta}_v^* M'$  in  $m$  steps, then  $M_m = M'$  by corollary 115.2, that is impossible because  $M_m \hat{\beta}_v M_{m+1}$  and so  $M_m$  is not  $\hat{\beta}_v$ -normal. □

The  $\hat{\beta}_v$ -reduction is not necessarily normalizing: if  $M = (\lambda x(x^1)x^1)^!(\lambda x(x^1)x^1)^!$  then  $M \hat{\beta}_v M$  and there is only one  $\hat{\beta}_v$ -redex in  $M$ , so  $M$  is a not  $\hat{\beta}_v$ -normalizable (closed) term. Moreover, the fact that a term is strongly  $\hat{\beta}_v$ -normalizable does not imply that it is  $\beta_v$ -normalizable: for instance, if  $M$  is as above then  $(\lambda z M)^!$  is a  $\hat{\beta}_v$ -normal form but  $(\lambda z M)^! \beta_v (\lambda z M)^!$ .

**Definition 116** (Leftmost and rightmost(-outermost) reduction). *We define two binary relations on  $\Lambda_t$ :*

- the weak leftmost(-outermost)  $\beta_v$ -reduction, denoted by  $\hat{\beta}_{vl}$ , whose rules are:

$$\frac{}{(\lambda x M)^! V^! \hat{\beta}_{vl} M[V/x]} \beta \quad \frac{M \hat{\beta}_{vl} M'}{MN \hat{\beta}_{vl} M'N} @_l$$

$$\frac{N \hat{\beta}_{vl} N'}{V^! N \hat{\beta}_{vl} V^! N'} @_r$$

- the weak rightmost(-outermost)  $\beta_v$ -reduction, denoted by  $\hat{\beta}_{vr}$ , whose rules are:

$$\frac{\frac{\frac{}{(\lambda x M)^! V^! \hat{\beta}_{vr} M[V/x]}{\beta} \quad \frac{M \hat{\beta}_{vr} M'}{MV^! \hat{\beta}_{vr} M'V^!} @_{lv}}{\frac{N \hat{\beta}_{vr} N'}{MN \hat{\beta}_{vr} MN'} @_r}}{\beta} \quad \frac{M \hat{\beta}_{vr} M'}{MV^! \hat{\beta}_{vr} M'V^!} @_{lv}}{\frac{N \hat{\beta}_{vr} N'}{MN \hat{\beta}_{vr} MN'} @_r} @_r$$

**Remark 117.** Clearly,  $\hat{\beta}_{vl}, \hat{\beta}_{vr} \subseteq \hat{\beta}_v$  (the proof is by induction on the length of the derivation of  $M \hat{\beta}_{vl} M'$  or  $M \hat{\beta}_{vr} M'$  respectively). As a consequence, every  $\hat{\beta}_v$ -normal form is  $\hat{\beta}_{vl}$ - and  $\hat{\beta}_{vr}$ -normal (in particular, for every value  $V$ ,  $V^!$  is  $\hat{\beta}_{vl}$ - and  $\hat{\beta}_{vr}$ -normal, see remark 111.4). The converse fails to hold; for instance, if  $I = (\lambda x x^!)$ ,  $x_i$  is a variable and  $V_i$  is a value for  $i \in \{1, 2\}$ , then

$$M = ((x_1^!)V_1^!)((II)(I)I)(x_2^!)V_2^!$$

is a  $\hat{\beta}_{vl}$ - and  $\hat{\beta}_{vr}$ -normal form but  $M \hat{\beta}_v^+ ((x_1^!)V_1^!)(I)(x_2^!)V_2^!$ .

We can characterize terms which are not  $\hat{\beta}_{vl}$ - or  $\hat{\beta}_{vr}$ -normal (see also remark 100.2).

**Theorem 118.** *Let  $M, M'$  be terms:  $M \hat{\beta}_{vl}$  (resp.  $\hat{\beta}_{vr}$ )  $M'$  iff there exist  $\ell \in \mathbb{N}$ , values  $V_0, \dots, V_\ell$  and terms  $L_0, \dots, L_\ell, L'_0, \dots, L'_\ell$  such that:*

- $(\lambda x N)^! V_0^! L_{01} \cdots L_{0k_0}$  (resp.  $L_0 = (L_{01}) \cdots (L_{0k_0})(\lambda x N)^! V_0^!$ ) and  $(N[V_0/x])L_{01} \cdots L_{0k_0}$  (resp.  $L'_0 = (L_{01}) \cdots (L_{0k_0})N[V_0/x]$ ) for some  $k_0 \in \mathbb{N}$  and terms  $N, L_{01}, \dots, L_{0k_0}$ ;
- for every  $1 \leq i \leq \ell$ , we have  $(V_i^!)L_{i-1}L_{i1} \cdots L_{ik_i}$  (resp.  $L_i = (L_{i1}) \cdots (L_{ik_i})(L_{i-1})V_i^!$ ) and  $(V_i^!)L'_{i-1}L_{i1} \cdots L_{ik_i}$  (resp.  $L'_i = (L_{i1}) \cdots (L_{ik_i})(L'_{i-1})V_i^!$ ) for some  $k_i \in \mathbb{N}$  and terms  $L_{i1}, \dots, L_{ik_i}$ ;
- $M = L_\ell$  and  $M' = L'_\ell$ .

Furthermore, both of these decompositions, if any, are unique.

**PROOF.** We prove the statement about  $\hat{\beta}_{vr}$ , the proof for the  $\hat{\beta}_{vl}$  case is perfectly symmetric.

$\Leftarrow$ : Proof by induction on  $\ell \in \mathbb{N}$ .

If  $\ell = 0$ , then  $M = (L_{01}) \cdots (L_{0k_0})(\lambda x N)^! V_0^!$  and  $M' = (L_{01}) \cdots (L_{0k_0})N[V_0/x]$  for some  $k_0 \in \mathbb{N}$ , value  $V_0$  and terms  $N, L_{01}, \dots, L_{0k_0}$ ; hence  $M \hat{\beta}_{vr} M'$  by applying the  $\beta$ -rule and  $k_0$  times the  $@_r$ -rule.

If  $\ell > 0$ , then  $M = (L_{\ell 1}) \cdots (L_{\ell k_\ell})(L_{\ell-1})V_\ell^!$  and  $M' = (L_{\ell 1}) \cdots (L_{\ell k_\ell})(L'_{\ell-1})V_\ell^!$  for some  $k_\ell \in \mathbb{N}$ , value  $V_\ell$  and terms  $L_{\ell 1}, \dots, L_{\ell k_\ell}$ ; by induction hypothesis,  $L_{\ell-1} \hat{\beta}_{vr} L'_{\ell-1}$ , so  $M \hat{\beta}_{vr} M'$  by applying the  $@_{vr}$ -rule and  $k_\ell$  times the  $@_r$ -rule.

$\Rightarrow$ : The uniqueness is obvious. The proof for the existence is by induction on the length of the derivation of  $M \hat{\beta}_v M'$ . Let us consider the last rule.

If the last rule is  $\beta$ , then  $M = (\lambda x N)^! V^!$  and  $M' = N[V/x]$  for some term  $N$  and value  $V$ , so we conclude by taking  $\ell = 0 = k_0$ .

If the last rule is  $@_{lv}$ , then  $M = NV^!$  and  $M' = N'V^!$  for some value  $V$  and terms  $N, N'$  such that  $N \hat{\beta}_{vr} N'$ ; by induction hypothesis, there exist  $\ell \in \mathbb{N}$ , terms  $L_0, \dots, L_\ell, L'_0, \dots, L'_\ell$  and values  $V_0, \dots, V_\ell$  such that:

- $L_0 = (L_{01}) \cdots (L_{0k_0})(\lambda x L)^! V_0^!$  and  $L'_0 = (L_{01}) \cdots (L_{0k_0}) L[V_0/x]$  for some  $k_0 \in \mathbb{N}$  and terms  $L, L_{01}, \dots, L_{0k_0}$ ;
- for every  $1 \leq i \leq \ell$ ,  $L_i = (L_{i1}) \cdots (L_{ik_i})(L_{i-1})V_i^!$  and  $L'_i = (L_{i1}) \cdots (L_{ik_i})(L'_{i-1})V_i^!$  for some  $k_i \in \mathbb{N}$  and terms  $L_{i1}, \dots, L_{ik_i}$ ;
- $N = L_\ell$  and  $N' = L'_\ell$ .

We conclude by taking  $L_{\ell+1} = M$ ,  $L'_{\ell+1} = M'$  and  $k_{\ell+1} = 0$ .

If the last rule is  $@_r$ , then  $M = N_1 N_2$  and  $M' = N_1 N'_2$  for some terms  $N_1, N_2, N'_2$  such that  $N_2 \hat{\beta}_{vr} N'_2$ ; by induction hypothesis, there exist  $\ell \in \mathbb{N}$ , terms  $L_0, \dots, L_\ell, L'_0, \dots, L'_\ell$  and values  $V_0, \dots, V_\ell$  such that:

- $L_0 = (L_{01}) \cdots (L_{0k_0})(\lambda x L)^! V_0^!$  and  $L'_0 = (L_{01}) \cdots (L_{0k_0}) L[V_0/x]$  for some  $k_0 \in \mathbb{N}$  and terms  $L, L_{01}, \dots, L_{0k_0}$ ;
- for every  $1 \leq i \leq \ell$ ,  $L_i = (L_{i1}) \cdots (L_{ik_i})(L_{i-1})V_i^!$  and  $L'_i = (L_{i1}) \cdots (L_{ik_i})(L'_{i-1})V_i^!$  for some  $k_i \in \mathbb{N}$  and terms  $L_{i1}, \dots, L_{ik_i}$ ;
- $N_2 = L_\ell$  and  $N'_2 = L'_\ell$ .

We can conclude by replacing in the sequence of  $L_i$ 's (resp.  $L'_i$ 's),  $L_\ell$  (resp.  $L'_\ell$ ) with  $M_\ell = (M_0) \cdots (M_{m_\ell})(L_{\ell-1})V_\ell^!$  (resp.  $M'_\ell = (M_0) \cdots (M_{m_\ell})(L'_{\ell-1})V_\ell^!$ ), with  $m_\ell = k_\ell + 1$ ,  $M_0 = N_1$  and  $M_j = L_{\ell j-1}$  for every  $1 \leq j \leq m_\ell$ , thus  $M = M_\ell$  and  $M' = M'_\ell$ .  $\square$

Theorem 118 says that in every term there exists at most one  $\hat{\beta}_v$ -redex that can be reduced by  $\hat{\beta}_{vl}$ - (resp.  $\hat{\beta}_{vr}$ -) reduction: theorem 118 might be seen also as a sort of definition of " $\hat{\beta}_{vl}$ - (resp.  $\hat{\beta}_{vr}$ -) redex".

**Corollary 119.**

1. *There are no critical pairs for the  $\hat{\beta}_{vl}$  (resp.  $\hat{\beta}_{vr}$ )-reduction: if  $M, N_1, N_2$  are terms such that  $M \hat{\beta}_{vl}$  (resp.  $\hat{\beta}_{vr}$ )  $N_1$  and  $M \hat{\beta}_{vl}$  (resp.  $\hat{\beta}_{vr}$ )  $N_2$ , then  $N_1 = N_2$ .*
2. *For every term  $M$ , it is  $\hat{\beta}_{vl}$  (resp.  $\hat{\beta}_{vr}$ )-normal iff either  $M = V^!$  for some value  $V$  or there exist  $\ell \in \mathbb{N}$ , a variable  $x$ , values  $V_0, \dots, V_\ell$  and terms  $L_0, \dots, L_\ell$  such that:*

- $L_0 = (x^!)V_0^!L_{01}\cdots L_{0k_0}$  (resp.  $L_0 = (L_{01})\cdots(L_{0k_0})(x^!)V_0^!$ ) for some  $k_0 \in \mathbb{N}$  and terms  $L_{01}, \dots, L_{0k_0}$ ;
  - for every  $1 \leq i \leq \ell$ , we have  $L_i = (V_i^!)L_{i-1}L_{i1}\cdots L_{ik_i}$  (resp.  $L_i = (L_{i1})\cdots(L_{ik_i})(L_{i-1})V_i^!$ ) for some  $k_i \in \mathbb{N}$  and terms  $L_{i1}, \dots, L_{ik_i}$ ;
  - $M = L_\ell$ .
3. Every closed term is a  $\hat{\beta}_v$ -normal form iff it is a  $\hat{\beta}_{v1}$ -normal form iff it is a  $\hat{\beta}_{vr}$ -normal form.
4. For every closed  $\hat{\beta}_v$ -normalizable term  $M$ , if  $M'$  is the  $\hat{\beta}_v$ -normal form of  $M$  then  $M \hat{\beta}_{v1}^* M'$  and  $M \hat{\beta}_{vr}^* M'$ .
5. For every term  $M$ , it is  $\hat{\beta}_v$ -normal iff either  $M = V^!$  for some value  $V$  or there exist  $\ell \in \mathbb{N}$ , a variable  $x$ , values  $V_0, \dots, V_\ell$  and terms  $L_0, \dots, L_\ell$  such that:
- $L_0 = (x^!)V_0^!L_{01}\cdots L_{0k_0}$  (resp.  $L_0 = (L_{01})\cdots(L_{0k_0})(x^!)V_0^!$ ) for some  $k_0 \in \mathbb{N}$  and  $\hat{\beta}_v$ -normal terms  $L_{01}, \dots, L_{0k_0}$ ;
  - for every  $1 \leq i \leq \ell$ , we have  $L_i = (V_i^!)L_{i-1}L_{i1}\cdots L_{ik_i}$  (resp.  $L_i = (L_{i1})\cdots(L_{ik_i})(L_{i-1})V_i^!$ ) for some  $k_i \in \mathbb{N}$  and  $\hat{\beta}_v$ -normal terms  $L_{i1}, \dots, L_{ik_i}$ ;
  - $M = L_\ell$ .

PROOF.

1. As every non- $\hat{\beta}_{v1}$  (resp.  $\hat{\beta}_{vr}$ )-normal term  $M$  can be written in a unique way in the forms of theorem 118, there is exactly one  $\hat{\beta}_v$ -redex in  $M$  that can be reduced by  $\hat{\beta}_{v1}$  (resp.  $\hat{\beta}_{vr}$ )-reduction.
2. It is an immediate consequence of theorem 118 and remark 100.2.
3. Every  $\hat{\beta}_v$ -normal form is obviously a  $\hat{\beta}_{v1}$ - (resp.  $\hat{\beta}_{vr}$ -)normal form (it is not necessary to suppose the term be closed, see remark 117).  
Conversely, let  $M$  be a closed  $\hat{\beta}_{v1}$ - (resp.  $\hat{\beta}_{vr}$ -)normal term: by corollary 119.2 and since  $M$  is closed, the only possibility is that  $M = V^!$  for some value  $V$ , so  $M$  is a  $\hat{\beta}_v$ -normal form by remark 111.4.
4. Proof by induction on the number  $n \in \mathbb{N}$  of steps of the  $\hat{\beta}_v$ -reduction from  $M$  to  $M'$  (this number is well-defined by corollary 115.2).  
If  $n = 0$  then  $M = M'$ , hence  $M \hat{\beta}_{v1}^* M'$  and  $M \hat{\beta}_{vr}^* M'$  (in 0 steps).  
If  $n > 0$  then  $M$  is not  $\hat{\beta}_v$ -normal. By corollary 119.3,  $M$  is neither a  $\hat{\beta}_{v1}$ - nor a  $\hat{\beta}_{vr}$ -normal form, hence there exist terms  $N_l$  and  $N_r$  such that  $M \hat{\beta}_{v1} N_l$  and  $M \hat{\beta}_{vr} N_r$ . As  $\hat{\beta}_{v1}, \hat{\beta}_{vr} \subseteq \hat{\beta}_v$ , both  $N_l \hat{\beta}_v^* M'$  and  $N_r \hat{\beta}_v^* M'$  in  $n - 1$  steps by corollaries 115.1-2. By induction hypothesis,  $N_l \hat{\beta}_{v1}^* M'$  and  $N_r \hat{\beta}_{vr}^* M'$ , thus  $M \hat{\beta}_{v1}^* M'$  and  $M \hat{\beta}_{vr}^* M'$  (in  $n$  steps).

5. If  $M$  is  $\hat{\beta}_v$ -normal then it is  $\hat{\beta}_{vl}$  (resp.  $\hat{\beta}_{vr}$ )-normal. By corollary 119.2, either  $M = V^!$  for some value  $V$  or there exist  $\ell \in \mathbb{N}$ , a variable  $x$ , values  $V_0, \dots, V_\ell$  and terms  $L_0, \dots, L_\ell$  such that:

- $L_0 = (x^!)V_0^!L_{01} \cdots L_{0k_0}$  (resp.  $L_0 = (L_{01}) \cdots (L_{0k_0})(x^!)V_0^!$ ) for some  $k_0 \in \mathbb{N}$  and terms  $L_{01}, \dots, L_{0k_0}$ ;
- for every  $1 \leq i \leq \ell$ , we have  $L_i = (V_i^!)L_{i-1}L_{i1} \cdots L_{ik_i}$  (resp.  $L_i = (L_{i1}) \cdots (L_{ik_i})(L_{i-1})V_i^!$ ) for some  $k_i \in \mathbb{N}$  and terms  $L_{i1}, \dots, L_{ik_i}$ ;
- $M = L_\ell$ ;

moreover, for every  $1 \leq i \leq \ell$  and  $1 \leq j \leq k_\ell$ ,  $L_{ij}$  is  $\hat{\beta}_v$ -normal (again since  $M$  is  $\hat{\beta}_v$ -normal) and so  $L_{ij}$  is  $\hat{\beta}_{vl}$  (resp.  $\hat{\beta}_{vr}$ )-normal by remark 117.

Conversely, let  $M$  be a term. If  $M = V^!$  for some value  $V$  then  $M$  is  $\hat{\beta}_v$ -normal (see remark 111.4).

If there exist  $\ell \in \mathbb{N}$ , a variable  $x$ , values  $V_0, \dots, V_\ell$  and terms  $L_0, \dots, L_\ell$  such that:

- $L_0 = (x^!)V_0^!L_{01} \cdots L_{0k_0}$  (resp.  $L_0 = (L_{01}) \cdots (L_{0k_0})(x^!)V_0^!$ ) for some  $k_0 \in \mathbb{N}$  and  $\hat{\beta}_v$ -normal terms  $L_{01}, \dots, L_{0k_0}$ ,
- for every  $1 \leq i \leq \ell$ , we have  $L_i = (V_i^!)L_{i-1}L_{i1} \cdots L_{ik_i}$  (resp.  $L_i = (L_{i1}) \cdots (L_{ik_i})(L_{i-1})V_i^!$ ) for some  $k_i \in \mathbb{N}$  and  $\hat{\beta}_v$ -normal terms  $L_{i1}, \dots, L_{ik_i}$ ,
- $M = L_\ell$ ,

then we show by induction on  $\ell \in \mathbb{N}$  that  $M$  is  $\hat{\beta}_v$ -normal and  $L_i$  is an application for every  $1 \leq i \leq \ell$ . If  $\ell = 0$  then  $M = L_0 = (x^!)V_0^!L_{01} \cdots L_{0k_0}$  which is an application (for any  $k_0 \in \mathbb{N}$ ) and a  $\hat{\beta}_v$ -normal form since  $L_{0j}$  is so for every  $1 \leq j \leq k_0$  by hypothesis. If  $\ell > 0$ , then  $M = L_\ell = (V_\ell^!)L_{\ell-1}L_{\ell 1} \cdots L_{\ell k_\ell}$  (resp.  $M = L_\ell = (L_{\ell 1}) \cdots (L_{\ell k_\ell})(L_{\ell-1})V_\ell^!$ ); by induction hypothesis,  $L_{\ell-1}$  is a  $\hat{\beta}_v$ -normal application, hence  $(V_\ell^!)L_{\ell-1}$  (resp.  $(L_{\ell-1})V_\ell^!$ ) is a  $\hat{\beta}_v$ -normal application; thus  $M$  is a  $\hat{\beta}_v$ -normal (since  $L_{\ell j}$  is so for every  $1 \leq j \leq k_\ell$  by hypothesis) application.  $\square$

In other words, according to corollary 119.1, the  $\hat{\beta}_{vl}$ - (resp.  $\hat{\beta}_{vr}$ -) reduction is “strongly deterministic” i.e. it is a partial map from  $\Lambda_t$  to  $\Lambda_t$ : any term  $M$  has at most one  $\hat{\beta}_{vl}$ - (resp.  $\hat{\beta}_{vr}$ -) redex, if any it is the “leftmost- (resp. rightmost-) outermost”  $\hat{\beta}_v$ -redex in  $M$  and there exists a unique term  $M'$  such that  $M \hat{\beta}_{vl}$  (resp.  $\hat{\beta}_{vr}$ )  $M'$ , otherwise if  $M$  is closed then it is  $\hat{\beta}_v$ -normal.

Corollary 119.2 provides a characterization of  $\hat{\beta}_{vl}$ - and  $\hat{\beta}_{vr}$ -normal forms. Corollary 119.5 claims that a term is  $\hat{\beta}_v$ -normal iff it is “hereditarily”  $\hat{\beta}_{vl}$ -normal iff it is “hereditarily”  $\hat{\beta}_{vr}$ -normal. These characterizations are more comprehensible by decomposing terms as binary trees (see §3.4).

The equivalences stated by corollary 119.3 have to be read together with the characterization given by proposition 113.

Corollary 119.4 provides two perfectly symmetric “ $\hat{\beta}_v$ -normalizing strategies”, which can be used for any  $\hat{\beta}_v$ -normalizable closed term.

Note that the hypothesis that the term is closed is necessary in corollaries 119.3-4: a term with some free variable might have a  $\hat{\beta}_v$ -redex without having neither “ $\hat{\beta}_{vl}$ ” nor “ $\hat{\beta}_{vr}$ -redex”, see for example the term  $M$  in remark 117, which is a  $\hat{\beta}_{vl}$ - and  $\hat{\beta}_{vr}$ -normal form but not a  $\hat{\beta}_v$ -normal form, so its  $\hat{\beta}_v$ -normal form cannot be reached by either  $\hat{\beta}_{vl}$ - or  $\hat{\beta}_{vr}$ -reduction.

We introduce now a “ $\hat{\beta}_v$ -normalization strategy” for which it is not necessary to assume that the term is closed.

**Definition 120** ( $\hat{\beta}_{vt}$ -reduction). *We define a relation  $\hat{\beta}_{vt} \subseteq \Lambda_t \times \Lambda_t$ , called turbo weak  $\hat{\beta}_v$ -reduction or  $\hat{\beta}_{vt}$ -reduction, by the following rules:*

$$\frac{}{(\lambda x M)!V! \hat{\beta}_{vt} M[V/x]} \beta \quad \frac{M \hat{\beta}_{vt} M' \quad N \hat{\beta}_{vt} N'}{MN \hat{\beta}_{vt} M'N'} @$$

$$\frac{M \hat{\beta}_{vt} M' \quad N \text{ is } \hat{\beta}_v\text{-normal}}{MN \hat{\beta}_{vt} M'N} @_{in}$$

$$\frac{N \hat{\beta}_{vt} N' \quad M \text{ is } \hat{\beta}_v\text{-normal}}{MN \hat{\beta}_{vt} MN'} @_m$$

*A term  $M$  is a  $\hat{\beta}_{vt}$ -normal form or is  $\hat{\beta}_{vt}$ -normal if there is no term  $M'$  such that  $M \hat{\beta}_{vt} M'$ .*

The following proposition clarifies the intuitive meaning of the  $\hat{\beta}_{vt}$ -reduction.

**Proposition 121.** *Let  $M, M'$  be terms:*

- $M \hat{\beta}_{vt} M'$  iff  $M$  contains at least one  $\hat{\beta}_v$ -redex and  $M'$  is obtained from  $M$  by replacing all the  $\hat{\beta}_v$ -redexes in  $M$  with their contractums;
- if  $M \hat{\beta}_{vt} M'$  then  $M \hat{\beta}_v^+ M'$  in  $n$  steps, where  $n$  is the number of  $\hat{\beta}_v$ -redexes in  $M$ .

PROOF.

$\Rightarrow$ : Proof by induction on the length of the derivation of  $M \hat{\beta}_{vt} M'$ . Let us consider the last rule of this derivation.

If it is the  $\beta$ -rule, then  $M = (\lambda x N)!V!$  and  $M' = N[V/x]$  for some term  $N$  and value  $V$ , so  $M$  is the only  $\hat{\beta}_v$ -redex in  $M$ ,  $M'$  is its contractum and  $M \hat{\beta}_v M'$  (in one step) by the  $\beta$ -rule.

If it is the @-rule, then  $M = M_1 M_2$  and  $M' = M'_1 M'_2$  for some terms  $M_1, M_2, M'_1, M'_2$  with  $M_i \hat{\beta}_{vt} M'_i$  for  $i \in \{1, 2\}$ ; by induction hypothesis,

$M_i$  contains a  $\hat{\beta}_v$ -redex,  $M'_i$  is obtained from  $M_i$  by replacing all the  $\hat{\beta}_v$ -redexes in  $M_i$  with their contractums and  $M_i \hat{\beta}_v^+ M'_i$  in  $n_i$  steps where  $n_i$  is the number of  $\hat{\beta}_v$ -redexes in  $M_i$ , for  $i \in \{1, 2\}$ .  $M$  is not a  $\hat{\beta}_v$ -redex, otherwise  $M_2 = V^!$  for some value  $V$  that is impossible by remark 103.1 since  $M_2$  contains a  $\hat{\beta}_v$ -redex. Hence  $M'$  is obtained from  $M$  by replacing all the  $\hat{\beta}_v$ -redexes in  $M$  with their contractums, moreover  $(M_1)M_2 \hat{\beta}_v^+ (M'_1)M_2$  in  $n_1$  steps and  $(M'_1)M_2 \hat{\beta}_v^+ (M'_1)M'_2$  in  $n_2$  steps, thus  $M \hat{\beta}_v^+ M'$  in  $n_1 + n_2$  steps, where  $n_1 + n_2$  is the number of  $\hat{\beta}_v$ -redexes in  $M$ .

If it is the  $@_{in}$ -rule, then  $M = M_1M_2$  and  $M' = M'_1M_2$  for some terms  $M_1, M_2, M'_1$  where  $M_1 \hat{\beta}_{vt} M'_1$  and  $M_2$  is  $\hat{\beta}_v$ -normal; by induction hypothesis,  $M_1$  contains a  $\hat{\beta}_v$ -redex,  $M'_1$  is obtained from  $M_1$  by replacing all the  $\hat{\beta}_v$ -redexes in  $M_1$  with their contractums and  $M_1 \hat{\beta}_v^+ M'_1$  in  $n$  steps where  $n$  is the number of  $\hat{\beta}_v$ -redexes in  $M_1$ .  $M_2$  contains no  $\hat{\beta}_v$ -redexes (by remark 111.1) and  $M$  is not a  $\hat{\beta}_v$ -redex (otherwise  $M_1 = (\lambda xN)^!$  for some term  $N$  that is impossible by remark 103.1 since  $M_1$  contains a  $\hat{\beta}_v$ -redex). Hence  $M'$  is obtained from  $M$  by replacing all the  $\hat{\beta}_v$ -redexes in  $M$  with their contractums, moreover  $M \hat{\beta}_v^+ M'$  in  $n$  steps, where  $n$  is the number of  $\hat{\beta}_v$ -redexes in  $M$ .

If it is the  $@_{rn}$ -rule, the proof is analogous to the previous case.

$\Leftarrow$ : Proof by induction on the term  $M$ . As  $M$  contains a  $\hat{\beta}_v$ -redex,  $M = M_1M_2$  for some terms  $M_1$  and  $M_2$  by remark 103.1.

If  $M_1$  and  $M_2$  contain no  $\hat{\beta}_v$ -redexes then  $M$  is the only  $\hat{\beta}_v$ -redex in  $M$ , so  $M_1 = (\lambda xN)^!$ ,  $M_2 = V^!$  and  $M' = N[V/x]$  (since  $M'$  is obtained by  $M$  by replacing all the  $\hat{\beta}_v$ -redexes in  $M$  with their contractums) for some term  $N$  and value  $V$ , thus  $M \hat{\beta}_{vt} M'$  by the  $\beta$ -rule.

If  $M_1$  and  $M_2$  contain a  $\hat{\beta}_v$ -redex, then by induction hypothesis  $M_i \hat{\beta}_{vt} M'_i$  where  $M'_i$  is obtained from  $M_i$  by replacing all the  $\hat{\beta}_v$ -redexes in  $M_i$  with their contractums, for  $i \in \{1, 2\}$ . Hence  $M' = M'_1M'_2$  since  $M$  is not a  $\hat{\beta}_v$ -redex (otherwise  $M_2 = V^!$  for some value  $V$  that is impossible by remark 103.1 since  $M_2$  contains a  $\hat{\beta}_v$ -redex). Thus  $M \hat{\beta}_{vt} M'$  by the  $@$ -rule.

If  $M_1$  contains a  $\hat{\beta}_v$ -redex and  $M_2$  does not, then  $M_2$  is  $\hat{\beta}_v$ -normal (by remark 103.1) and  $M_1 \hat{\beta}_{vt} M'_1$  where  $M'_1$  is obtained from  $M_1$  by replacing all the  $\hat{\beta}_v$ -redexes in  $M_1$  with their contractums (by induction hypothesis). Hence  $M' = M'_1M_2$  since  $M_2$  contains no  $\hat{\beta}_v$ -redexes (by remark 111.1) and  $M$  is not a  $\hat{\beta}_v$ -redex (otherwise  $M_1 = (\lambda xN)^!$  for some term  $N$  that is impossible by remark 103.1 since  $M_1$  contains a  $\hat{\beta}_v$ -redex). Thus  $M \hat{\beta}_{vt} M'$  by the  $@_{in}$ -rule.

If  $M_2$  contains a  $\hat{\beta}_v$ -redex and  $M_1$  does not, the proof is analogous to the previous case, in particular we conclude that  $M \hat{\beta}_{vt} M'$  by applying



the  $@_{rn}$ -rule. □

**Corollary 122.**

1. Every term is  $\hat{\beta}_v$ -normal iff it is  $\hat{\beta}_{vt}$ -normal.
2. Terms of the shape  $V^!$  for some value  $V$  are  $\hat{\beta}_{vt}$ -normal.
3. There are no critical pairs for the  $\hat{\beta}_{vt}$ -reduction: if  $M, N_1, N_2$  are terms such that  $M \hat{\beta}_{vt} N_1$  and  $M \hat{\beta}_{vt} N_2$ , then  $N_1 = N_2$ .

PROOF.

1. By proposition 121 and remark 111.1.
2. By corollary 122.1 ( $\Rightarrow$ ) and remark 111.4.
3. Immediate consequence of proposition 121. □

Corollary 122.3 says that the  $\hat{\beta}_{vt}$ -reduction is “strongly deterministic” (i.e. it is a partial map from  $\Lambda_t$  to  $\Lambda_t$ ): if a term  $M$  is not  $\hat{\beta}_v$ -normal, then there exists a unique term  $M'$  such that  $M \hat{\beta}_{vt} M'$ .

Obviously, the fact that  $M \hat{\beta}_{vt} M'$  does not entail that  $M'$  is  $\hat{\beta}_v$ -normal, since the  $\hat{\beta}_{vt}$ -reduction might create new  $\hat{\beta}_v$ -redexes. For instance  $(\lambda x(x^!)x^!)^!(\lambda x x^!)^! \hat{\beta}_{vt} (\lambda x x^!)^!(\lambda x x^!)^!$  which is not  $\hat{\beta}_v$ -normal.

**Theorem 123.** *For every  $\hat{\beta}_v$ -normalizable term  $M$ , if  $M'$  is the  $\hat{\beta}_v$ -normal form of  $M$  then  $M \hat{\beta}_{vt}^* M'$ .*

PROOF. By induction on the number  $m \in \mathbb{N}$  of steps of the  $\hat{\beta}_v$ -reduction  $M \hat{\beta}_v^* M'$ .

If  $m = 0$  then  $M = M'$ , therefore  $M \hat{\beta}_{vt}^* M'$  by reflexivity of  $\hat{\beta}_{vt}^*$ .

If  $m > 0$  then  $M$  is not  $\hat{\beta}_v$ -normal, thus  $M$  is not  $\hat{\beta}_{vt}$ -normal by corollary 122.1, hence there exists a term  $N$  such that  $M \hat{\beta}_{vt} N$ . By proposition 121,  $M \hat{\beta}_v^+ N$  in  $n > 0$  steps, with  $n \leq m$  by corollary 115.2. By corollary 115.1 there exists a term  $N'$  such that  $M' \hat{\beta}_v^* N'$  and  $N \hat{\beta}_v^* N'$ , so  $M' = N'$  since  $M'$  is  $\hat{\beta}_v$ -normal and thus  $N \hat{\beta}_{vt}^* M'$  in  $m - n < m$  by corollary 115.2. Therefore  $N \hat{\beta}_{vt}^* M'$  by induction hypothesis, hence  $M \hat{\beta}_{vt}^* M'$  by transitivity of  $\hat{\beta}_{vt}^*$ . □

Theorem 123 provides a “ $\hat{\beta}_v$ -normalizing strategy”, which can be used for any  $\hat{\beta}_v$ -normalizable (not necessarily closed) term.

### 3.1.3 Some problems with $\eta$ -reduction

**Definition 124** ( $\eta$ -reduction). *We define a relation  $\eta \subseteq (\Lambda_t \times \Lambda_t) \cup (\Lambda_v \times \Lambda_v)$ , called  $\eta$ -reduction, by the following rules:*

$$\frac{}{(\lambda x M x)! \eta M} \eta \quad \frac{M \eta M'}{MN \eta M'N} @_l \quad \frac{N \eta N'}{MN \eta MN'} @_r$$

$$\frac{M \eta M'}{\lambda x M \eta \lambda x M'} \lambda \quad \frac{V \eta W}{V! \eta W!} !$$

where in  $\eta$ -rule the variable  $x$  is not free in the term  $M$ .

**Remark 125.** Let  $M = (\lambda d y^!)^! (\lambda z ((\lambda x x^! x^!)^! (\lambda x x^! x^!)^!) z^!)^!$ : then  $M \eta (\lambda d y^!)^! ((\lambda x x^! x^!)^! (\lambda x x^! x^!)^!)$   $M \hat{\beta}_v y^!$  where  $y^!$  and  $(\lambda d y^!)^! ((\lambda x x^! x^!)^! (\lambda x x^! x^!)^!)$  are  $\beta_v \eta$ -normal forms. Therefore, neither the  $\beta_v \eta$ -reduction nor the  $\hat{\beta}_v \eta$ -reduction are confluent.

### 3.2 A “completion” of $\beta_v$ -reduction

A solid theory of call-by-value  $\lambda$ -calculus requires an operational characterization of solvability, i.e. to find a strategy which computes the results of the represented functions. Following [PR04], a term  $t$  is CBV-solvable whenever there is an head context  $H$  s.t.  $H[t] \rightarrow_{\beta_v}^* I$  where  $I = \lambda x x$  and  $\rightarrow_{\beta_v}$  is the call-by-value  $\beta$ -reduction in Plotkin’s  $\lambda_{\beta_v}$ -calculus. An operational characterization has been provided in [PR99, PR04] but, unfortunately, it is obtained through call-by-name  $\beta$ -reduction, which is disappointing and not satisfying. If it is not possible to get an internal characterization, i.e. one which uses the rules of the calculus itself, then there is an inherent weakness in the rewriting rules of the calculus. For Plotkin’s call-by-value  $\lambda_{\beta_v}$ -calculus [Plo75] it is indeed the case, let us illustrate the point with an example. Let  $\Delta = \lambda x (x)x$ . There is no head context sending (via  $\beta_v$ -reduction) the following term to the identity:

$$t = ((\lambda y \Delta)(x)z)\Delta$$

and – as a consequence –  $t$  should be unsolvable and divergent in a good call-by-value calculus, while it is in  $\beta_v$ -normal form. The weakness of  $\beta_v$ -reduction is a fact widely recognized and accepted, indeed there have been many proposals of alternative call-by-value  $\lambda$ -calculi, see for instance [Mog89, Hof95, DL07, HZ09, AP12]. All these different versions of call-by-value  $\lambda$ -calculi extend the syntax of  $\lambda$ -calculus with an explicit substitution constructor  $t\{u/x\}$  (which is equivalent to use **let**...**in** expressions) defined in the syntax, but these substitutions are just delayed, they are not propagated in a small-steps way.

In particular, Accattoli and Paolini introduced in [AP12] the value-substitution lambda-calculus, a simple call-by-value  $\lambda$ -calculus with explicit substitutions borrowing ideas from Herbelin and Zimmerman’s lambda-CBV calculus ([HZ09]) and from Accattoli and Kesner’s structural lambda-calculus ([AK10]), both with explicit substitutions. Interestingly, in this new setting, Accattoli and Paolini characterized solvable terms as those terms having normal form with respect to a suitable contextual closure of its (call-by-value) reduction rules, thus improving over the previous characterization.

We aim at showing that we can characterize CBV-solvable terms without using explicit substitutions, by only adding some simple reduction rules in our syntax. These supplementary rules are nothing but an orientation of the two orientable rules  $\sigma_1$  and  $\sigma_3$  generating the  $\sigma_v$ -equivalence (see section 5.2): they are a reformulation in our syntax without explicit substitutions of the  $let_{let}$ - and  $let_{app}$ -rules of the Herbelin’s and Zimmerman’s calculus (see [HZ09]).

**Definition 126** ( $\sigma$ - and  $\sigma'$ -reduction).  $\sigma_1$  is the contextual closure of the binary relation  $\rightarrow_{\sigma_1}$  on  $\Lambda_t$  defined by the following rule:

$$\frac{}{(\lambda x M)^! N L \rightarrow_{\sigma_1} (\lambda x M L)^! N} \sigma_1$$

where  $M, N, L$  are terms and  $x \notin \text{fv}(L)$ .

$\sigma_3$  is the contextual closure of the binary relation  $\rightarrow_{\sigma_3}$  on  $\Lambda_t$  defined by the following rule:

$$\frac{}{(V^!)((\lambda x L)^!) N \rightarrow_{\sigma_3} (\lambda x V^! L)^! N} \sigma_3$$

where  $N$  and  $L$  are terms,  $V$  is a value and  $x \notin \text{fv}(V)$ .

$\sigma'_3$  is the contextual closure of the binary relation  $\rightarrow_{\sigma'_3}$  on  $\Lambda_t$  defined by the following rule:

$$\frac{}{(M)((\lambda x L)^!) N \rightarrow_{\sigma'_3} (\lambda x M L)^! N} \sigma'_3$$

where  $M, N, L$  are terms and  $x \notin \text{fv}(M)$ .

The  $\sigma_v$ -reduction (resp.  $\sigma'_v$ -reduction) is  $\sigma_v = \sigma_1 \cup \sigma_3$  (resp.  $\sigma'_v = \sigma_1 \cup \sigma'_3$ ).

The variable condition on  $\sigma_1$ -,  $\sigma_3$ - and  $\sigma'_3$ -rules can be always fulfilled by  $\alpha$ -conversion.

The  $\sigma_3$ -rule is a weakened version of the  $\sigma'_3$ -rule, i.e. it is the  $\sigma'_3$ -rule limited to the case where  $M = V^!$  for some value  $V$ .

$\sigma_1$ - and  $\sigma'_3$ -rules above are just an orientation of respective rules in the definition of  $\sigma_v$ -equivalence. Note the left-right symmetry of  $\sigma_1$ - and  $\sigma'_3$ -rules: in  $\sigma_1$  (resp.  $\sigma'_3$ )-rule, the  $\sigma$ -redex is an application of a  $\beta$ -redex (resp. term) to a term (resp.  $\beta$ -redex). In remark 135, we will see a reason to like the  $\sigma_3$ -rule more than its generalization  $\sigma'_3$ .

**Remark 127.**

1. By remark 105, one has  $\sigma_v, \sigma'_v \subseteq (\Lambda_t \times \Lambda_t) \cup (\Lambda_v \times \Lambda_v)$ . Moreover,  $\sigma_v$  (resp.  $\sigma'_v$ ) is the contextual closure of the relations  $\rightarrow_{\sigma_3}$  (resp.  $\rightarrow_{\sigma'_3}$ ) and  $\rightarrow_{\sigma_1}$  (the proof is by straightforward induction on the derivations).
2. It is immediate to check that for every expression  $E$ , there exists an expression  $E'$  such that  $E \sigma_v E'$  (resp.  $E \sigma'_v E'$ ) iff  $E$  contains a  $\sigma_v$  (resp.  $\sigma'_v$ )-redex. Therefore, an expression is  $\sigma_v$  (resp.  $\sigma'_v$ )-normal iff it contains no  $\sigma_v$  (resp.  $\sigma'_v$ )-redex.

3. It is easy to verify that for all expressions  $E, E'$ ,  $E \sigma_v E'$  (resp.  $E \sigma'_v E'$ ) iff  $E'$  is obtained from  $E$  by replacing exactly one  $\sigma_v$  (resp.  $\sigma'_v$ )-redex in  $E$  with its contractum.
4. Clearly,  $\sigma_v \subseteq \sigma'_v$ . The converse does not hold, for instance take  $M = (z_1^! z_2^!)$ ,  $N = (y_1^! y_2^!)$  and  $L = (M)((\lambda x_1 x_2^!)^! N)$  where  $x_1, x_2, y_1, y_2, z_1, z_2$  are pairwise distinct variables: then  $L$  is  $\sigma_v$ -normal but  $L \sigma'_3 (\lambda x_1 (M) x_2^!)^! N$ .

We can merge the  $\sigma_v$ - and  $\sigma'_v$ -reduction into the  $\beta_v$ - and  $\hat{\beta}_v$ -reduction, in order to get a sort of “completion” of the  $\beta_v$ - and  $\hat{\beta}_v$ -reduction.

**Definition 128.** We set  $\beta_{v\sigma} = \beta_v \cup \sigma_v$  (resp.  $\beta_{v\sigma'} = \beta_v \cup \sigma'_v$ ), called  $\beta_{v\sigma}$  (resp.  $\beta_{v\sigma'}$ )-reduction, and  $\hat{\beta}_{v\sigma} = \hat{\beta}_v \cup \sigma_v$  (resp.  $\hat{\beta}_{v\sigma'} = \hat{\beta}_v \cup \sigma'_v$ ), called  $\hat{\beta}_{v\sigma}$  (resp.  $\hat{\beta}_{v\sigma'}$ )-reduction.

Intuitively,  $\sigma_v$ -reduction might enable a  $\beta_v$ -redex in an expression  $E$  which is hidden by the inessential sequential structure of  $E$ . For instance, if  $N = (z^!)z^!$  and  $M = (\lambda y (\lambda x z_0^!)^! x^!)^! N$  where  $x, y, z, z_0$  are pairwise distinct variables, then  $M_1 = (\lambda y (\lambda x z_0^!)^!)^! N x^!$  and  $M_2 = (\lambda x z_0^!)^! ((\lambda y x^!)^!) N$  are  $\beta_v$ -normal, but  $M_1 \sigma_v M$  (by the  $\sigma_1$ -rule) and  $M_2 \sigma_v M$  (by the  $\sigma_3$ -rule), where  $M$  is not  $\beta_v$ -normal.

**Remark 129.**

1. By remark 105, one has  $\beta_{v\sigma}, \hat{\beta}_{v\sigma}, \beta_{v\sigma'}, \hat{\beta}_{v\sigma'} \subseteq (\Lambda_t \times \Lambda_t) \cup (\Lambda_v \times \Lambda_v)$ .
2. It is immediate to check that for every expression  $E$  (resp. term  $M$ ), there exists an expression  $E'$  (resp. a term  $M'$ ) such that  $E \beta_{v\sigma} E'$  (resp.  $M \hat{\beta}_{v\sigma} M'$ ) iff  $E$  (resp.  $M$ ) contains a  $\beta_{v\sigma}$  (resp.  $\hat{\beta}_{v\sigma}$ )-redex. Therefore, an expression (resp. a term) is  $\beta_{v\sigma}$  (resp.  $\hat{\beta}_{v\sigma}$ )-normal iff it contains no  $\beta_{v\sigma}$  (resp.  $\hat{\beta}_{v\sigma}$ )-redex. Analogous considerations hold for  $\beta_{v\sigma'}$  (resp.  $\hat{\beta}_{v\sigma'}$ )-reduction.
3. It is easy to verify that for all expressions  $E, E'$  (resp. terms  $M, M'$ ),  $E \beta_{v\sigma} E'$  (resp.  $M \hat{\beta}_{v\sigma} M'$ ) iff  $E'$  (resp.  $M'$ ) is obtained from  $E$  (resp.  $M$ ) by replacing exactly one  $\beta_{v\sigma}$  (resp.  $\hat{\beta}_{v\sigma}$ )-redex in  $E$  (resp.  $M$ ) with its contractum. Analogous considerations hold for  $\beta_{v\sigma'}$  (resp.  $\hat{\beta}_{v\sigma'}$ )-reduction.
4. Clearly,  $\hat{\beta}_{v\sigma} \subseteq \beta_{v\sigma} \subseteq \beta_{v\sigma'}$  and  $\hat{\beta}_{v\sigma'} \subseteq \hat{\beta}_{v\sigma} \subseteq \beta_{v\sigma'}$ . The converses do not hold.

We prove now a confluence property for  $\beta_{v\sigma}$  and  $\hat{\beta}_{v\sigma}$ . For this purpose, we use a commutation property of  $\beta_v$ - and  $\sigma_v$ -reductions and the strong normalization of  $\sigma'_v$ .

**Definition 130.** With every expression  $E$  are associated two measures  $\text{size}'(E), \#_w(E) \in \mathbb{N}$ , defined by induction on  $E$  as follows:

- $\text{size}'(x) = 2$ ;
- $\text{size}'(\lambda x M) = \text{size}'(M) + 1$ ;
- $\text{size}'(V^!) = \text{size}'(V)$ ;
- $\text{size}'(MN) = \text{size}'(M) + \text{size}'(N)$ .
- $\#_w(x) = 1$ ;
- $\#_w(\lambda x M) = \#_w(M) + \text{size}'(M)$ ;
- $\#_w(V^!) = \#_w(V)$ ;
- $\#_w(MN) = \#_w(M) + \#_w(N) + 2\text{size}'(M)\text{size}'(N) - 1$ .

$\#_w(M)$  is the sum of the weights of the nodes in  $\mathcal{T}_@M$ , where the weight of a node  $n$  in  $\mathcal{T}_@M$  is the difference between  $\text{size}(M)$  and the number of  $\lambda$ -nodes above  $n$ .

WHY?

**Remark 131.**  $\text{size}'(E) \geq 2$  and  $\#_w(E) \geq 1$  for any expression  $E$ . The proof is by a straightforward induction on the expression  $E$ .

**Lemma 132.** Let  $E$  and  $E'$  be expressions. If  $E \sigma'_V E'$  then  $\#_w(E) > \#_w(E')$  and  $\text{size}'(E) = \text{size}'(E')$ .

PROOF. By induction on the length of the derivation of  $E \sigma'_V E'$ . Let us consider the last rule of this derivation.

If it is the  $\sigma_1$ -rule then  $E = (\lambda x M)^! N L$  and  $E' = (\lambda x M L)^! N$  for some terms  $N, M$  and  $L$ , so

$$\begin{aligned} \#_w(E) &= \#_w(M) + \#_w(L) + \#_w(N) + \text{size}'(M) + 2\text{size}'(N) + 2\text{size}'(L) + \\ &\quad 2\text{size}'(M)\text{size}'(N) + 2\text{size}'(M)\text{size}'(L) + 2\text{size}'(L)\text{size}'(N) - 2 \\ \#_w(E') &= \#_w(M) + \#_w(L) + \#_w(N) + \text{size}'(M) + 2\text{size}'(N) + \text{size}'(L) + \\ &\quad 2\text{size}'(M)\text{size}'(N) + 2\text{size}'(M)\text{size}'(L) + 2\text{size}'(L)\text{size}'(N) - 2 = \#_w(E) - \text{size}'(L) \end{aligned}$$

hence  $\#_w(E) > \#_w(E')$  by remark 131. Moreover,  $\text{size}'(E) = \text{size}'(M) + \text{size}'(L) + \text{size}'(N) + 1 = \text{size}'(E')$ .

If it is the  $\sigma_3$ -rule then  $E = (M)((\lambda x L)^!)N$  and  $E' = (\lambda x M L)^! N$  for some terms  $N, M$  and  $L$ , so

$$\begin{aligned} \#_w(E) &= \#_w(M) + \#_w(L) + \#_w(N) + 2\text{size}(M) + 2\text{size}(N) + \text{size}(L) + \\ &\quad 2\text{size}(M)\text{size}(N) + 2\text{size}(M)\text{size}(L) + 2\text{size}(L)\text{size}(N) - 2 \\ \#_w(E') &= \#_w(M) + \#_w(L) + \#_w(N) + \text{size}'(M) + 2\text{size}'(N) + \text{size}'(L) + \\ &\quad 2\text{size}'(M)\text{size}'(N) + 2\text{size}'(M)\text{size}'(L) + 2\text{size}'(L)\text{size}'(N) - 2 = \#_w(E) - \text{size}'(M) \end{aligned}$$

hence  $\#_w(E) > \#_w(E')$  by remark 131. Moreover,  $\text{size}'(E) = \text{size}'(M) + \text{size}'(L) + \text{size}'(N) + 1 = \text{size}'(E')$ .

If it is the  $\lambda$ -rule then  $E = \lambda x M$  and  $E' = \lambda x M'$  for some terms  $M$  and  $M'$  such that  $M \sigma'_v M'$ , thus  $\#_w(E) = \#_w(M) + \text{size}'(M)$  and  $\#_w(E') = \#_w(M') + \text{size}'(M')$ ; by induction hypothesis, we have  $\#_w(M) > \#_w(M')$  and  $\text{size}'(M) = \text{size}'(M')$ , therefore  $\#_w(E) > \#_w(M') + \text{size}'(M) = \#_w(E')$  and  $\text{size}'(E) = \text{size}'(M) + 1 = \text{size}'(E')$ .

If it is the  $!$ -rule then  $E = V^!$  and  $E' = V'^!$  for some values  $V$  and  $V'$  such that  $V \sigma'_v V'$ , thus  $\#_w(E) = \#_w(V)$  and  $\#_w(E') = \#_w(V')$ ; by induction hypothesis, we have  $\#_w(V) > \#_w(V')$  and  $\text{size}'(V) = \text{size}'(V')$ , therefore  $\#_w(E) > \#_w(E')$  and  $\text{size}'(E) = \text{size}'(V) = \text{size}'(E')$ .

If it is the  $@_l$  (resp.  $@_r$ )-rule then  $E = MN$  and  $E' = M'N$  (resp.  $E = MN'$ ) for some terms  $M, N$  and  $M'$  (resp.  $N'$ ) such that  $M \sigma'_v M'$  (resp.  $N \sigma'_v N'$ ), thus  $\#_w(E) = \#_w(M) + \#_w(N) + 2\text{size}'(M)\text{size}'(N) - 1$  and  $\#_w(E') = \#_w(M') + \#_w(N) + 2\text{size}'(M')\text{size}'(N) - 1$  (resp.  $\#_w(E') = \#_w(M) + \#_w(N') + 2\text{size}'(M)\text{size}'(N') - 1$ ); by induction hypothesis, we have  $\#_w(M) > \#_w(M')$  (resp.  $\#_w(N) > \#_w(N')$ ) and  $\text{size}'(M) = \text{size}'(M')$  (resp.  $\text{size}'(N) = \text{size}'(N')$ ), therefore  $\#_w(E) > \#_w(M) + \#_w(N) + 2\text{size}'(M)\text{size}'(N) - 1 = \#_w(E')$  (resp.  $\#_w(E) > \#_w(M) + \#_w(N) + 2\text{size}'(M)\text{size}'(N') - 1 = \#_w(E')$ ) and  $\text{size}'(E) = \text{size}'(M) + \text{size}'(N) = \text{size}'(E')$ .  $\square$

**Proposition 133.**  $\sigma'_v$  (and in particular  $\sigma_v$ ) is strongly normalizing.

PROOF. It is an immediate consequence of the previous lemma.  $\square$

**Lemma 134.**  $\sigma_v$  is locally confluent.

PROOF. By induction on the expression  $E$  such that  $E \sigma_v E_i$  for  $i \in \{1, 2\}$ . The only interesting case are:

- if  $E = ((\lambda x M)^!((\lambda y L)^!)N)L'$  with  $E \sigma_v (\lambda x M L')^!((\lambda y L)^!)N = E_1$  (by reducing the  $\sigma_1$ -redex  $E$ ) and  $E \sigma_v (\lambda y (\lambda x M)^!L)^!N L' = E_2$  (by reducing the  $\sigma_3$ -redex  $(\lambda x M)^!((\lambda y L)^!)N$  in  $E$ ), then  $E_2 \sigma_v (\lambda y (\lambda x M)^!L L')^!N \sigma_v (\lambda y (\lambda x M L')^!L)^!E'$  (by reducing twice a  $\sigma_1$ -redex) and  $E_1 \sigma_v E'$  (by reducing the  $\sigma_3$ -redex  $E_1$ );
- if  $E = (V^!)((\lambda x L)^!((\lambda x' L')^!)N)$  with  $E \sigma_v (V^!)(\lambda x' (\lambda x L)^!L')^!N = E_1$  (by reducing the  $\sigma_3$ -redex  $(\lambda x L)^!((\lambda x' L')^!)N$  in  $E$ ) and  $E \sigma_v (\lambda x V^!L)^!((\lambda x' L')^!)N = E_2$  (by reducing the  $\sigma_3$ -redex  $E$ ), then  $E_1 \sigma_v (\lambda x' (V)(\lambda x L)L')^!N \sigma_v (\lambda x' (\lambda x V^!L)^!L')^!N = E'$  (by reducing twice a  $\sigma_3$ -redex) and  $E_2 \sigma_v E'$  (by reducing the  $\sigma_3$ -redex  $E_2$ ).  $\square$

**Remark 135.**  $\sigma'_v$  and  $\beta_{v\sigma'}$  are not locally confluent and so neither confluent. For instance, take  $N_i = (z_i^!)z_i^!$  for  $i \in \{1, 2\}$  and  $M = ((\lambda x_1 y_1^!)^!N_1)((\lambda x_2 y_2^!)^!)N_2$  where  $x_1, x_2, y_1, y_2, z_1, z_2$  are pairwise distinct variables:  $M$  is  $\beta_v$ -normal and it contains no  $\sigma_3$ -redexes but  $M \sigma'_v (\lambda x_2 (\lambda x_1 y_1^!)^!N_1 y_2^!)^!N_2$  (because of the  $\sigma'_3$ -rule) which contains only a  $\sigma_1$ -redex, and  $M \sigma'_v (\lambda x_1 (y_1^!)(\lambda x_2 y_2^!)^!N_2)^!N_1$  (because of the  $\sigma_1$ -rule) which contains only a  $\sigma_3$ -redex, so  $M$  reduces to two different  $\sigma'_v$ -normal forms,  $(\lambda x_2 (\lambda x_1 y_1^!)^!N_1)^!N_2$  and  $(\lambda x_1 (\lambda x_2 y_1^!)^!N_2)^!N_1$ .

We conjecture that  $\sigma'$  is confluent modulo the equivalence relation on  $\Lambda_{\text{CBV}}$  generated by the following binary relation  $\sim_{\sigma_4}$  on  $\Lambda_t$  defined by:

$$(\lambda x_1(\lambda x_2 M)^! N_2)^! N_1 \sim_{\sigma_4} (\lambda x_2(\lambda x_1 M)^! N_1)^! N_2$$

where  $x_2 \notin \text{fv}(N_1)$  and  $x_1 \notin \text{fv}(N_2)$ .

**Proposition 136.**  $\sigma_v$  is confluent.

PROOF. By lemma 134, proposition 133 and Newman’s lemma.  $\square$

We recall a well-known result on term rewriting systems.

**Lemma 137** (Hindley–Rosen). *Let  $\rightarrow_1$  and  $\rightarrow_2$  be two binary relations on a set  $X$ . If they are both confluent and they commute, i.e. if  $t \rightarrow_1^* u_1$  and  $t \rightarrow_2^* u_2$  then there exists  $s$  such that  $u_1 \rightarrow_2^* s$  and  $u_2 \rightarrow_1^* s$ , then  $\rightarrow_1 \cup \rightarrow_2$  is confluent.*

PROOF. See proposition 3.3.5 in [Bar84].  $\square$

**Lemma 138.** *Let  $E$  and  $E'$  be expressions, let  $V$  and  $V'$  be values and let  $x$  be a variable:*

1. if  $V \sigma_v V'$  then  $E[V/x] \sigma_v^* E[V'/x]$ ;
2. if  $E \sigma_v E'$  then  $E[V/x] \sigma_v E'[V/x]$ .

PROOF.

1. By induction on the expression  $E$ .

If  $E = x$ , then  $E[V/x] = V$  and  $E[V'/x] = V'$ , so  $E[V/x] \sigma_v^* E[V'/x]$  by hypothesis.

If  $E = y \neq x$ , then  $E[V/x] = y = E[V'/x]$ , then  $E[V/x] \sigma_v^* E[V'/x]$  by reflexivity of  $\sigma_v^*$ .

If  $E = \lambda y M$  for some term  $M$ , then we can suppose without loss of generality that  $y \neq x$ , hence  $E[V/x] = \lambda y M[V/x]$  and  $E[V'/x] = \lambda y M[V'/x]$ ; by induction hypothesis,  $M[V/x] \sigma_v^* M[V'/x]$  and thus  $E[V/x] \sigma_v^* E[V'/x]$  since  $\sigma_v^*$  passes to context.

If  $E = W^!$  for some value  $W$ , then  $E[V/x] = (W[V/x])^!$  and  $E[V'/x] = (W[V'/x])^!$ ; by induction hypothesis,  $W[V/x] \sigma_v^* W[V'/x]$  and thus  $E[V/x] \sigma_v^* E[V'/x]$  since  $\sigma_v^*$  passes to context.

If  $E = MN$  for some terms  $M, N$ , then  $E[V/x] = M[V/x]N[V/x]$  and  $E[V'/x] = M[V'/x]N[V'/x]$ ;  $M[V/x] \sigma_v^* M[V'/x]$  and  $N[V/x] \sigma_v^* N[V'/x]$  by induction hypothesis, therefore  $E[V/x] \sigma_v^* M[V'/x]N[V'/x]$  since  $\sigma_v^*$  passes to context.

2. By induction on the length of the derivation of  $E \sigma_v E'$ . Let us consider the last rule of this derivation.

If it is the  $\sigma_1$ -rule, then  $E = (\lambda y M)^! N L$  and  $E' = (\lambda y M L)^! N$  with  $y \notin \text{fv}(L)$ ; we can suppose without loss of generality that  $y \notin \text{fv}(V) \cup \{x\}$ , hence  $E[V/x] = (\lambda y M[V/x])^! N[V/x] L[V/x]$  and  $E'[V/x] = (\lambda y M[V/x] L[V/x])^! N[V/x]$ , therefore  $E[V/x] \sigma_v E'[V/x]$  by the  $\sigma_1$ -rule, since  $y \notin (\text{fv}(L) \setminus \{x\}) \cup \text{fv}(V) = \text{fv}(L[V/x])$ .

If it is the  $\sigma_3$ -rule, then  $E = (W^!)((\lambda y L)^!) N$  and  $E' = (\lambda y W^! L)^! N$  with  $y \notin \text{fv}(W)$ ; we can suppose without loss of generality that  $y \notin \text{fv}(V) \cup \{x\}$ , so  $E[V/x] = (W[V/x])^!((\lambda y L[V/x])^!) N[V/x]$  and  $E'[V/x] = (\lambda y (W[V/x])^! L[V/x])^! N[V/x]$ , therefore  $E[V/x] \sigma_v E'[V/x]$  by the  $\sigma_3$ -rule, since  $y \notin (\text{fv}(W) \setminus \{x\}) \cup \text{fv}(V) = \text{fv}(W[V/x])$ .

If it is the  $\lambda$ -rule then  $E = \lambda y M$  and  $E' = \lambda y M'$  for some terms  $M$  and  $M'$  with  $M \sigma_v M'$ ; we can suppose without loss of generality that  $y \notin \text{fv}(V) \cup \{x\}$ , hence  $E[V/x] = \lambda y M[V/x]$  and  $E'[V/x] = \lambda y M'[V/x]$ ; by induction hypothesis,  $M[V/x] \sigma_v M'[V/x]$  and thus  $E[V/x] \sigma_v E'[V/x]$  since by the  $\lambda$ -rule.

If it is the  $!$ -rule then  $E = W^!$  and  $E' = W'^!$  for some values  $W$  such that  $W \sigma_v W'$ , so  $E[V/x] = (W[V/x])^!$  and  $E'[V/x] = (W'[V/x])^!$ ; by induction hypothesis,  $W[V/x] \sigma_v W'[V/x]$  and thus  $E[V/x] \sigma_v E'[V/x]$  by the  $!$ -rule.

If it is the  $@_l$  (resp.  $@_r$ ) then  $E = M N$  and  $E' = M' N$  (resp.  $E' = M N'$ ) for some terms  $M, N$  and  $M'$  (resp.  $N'$ ) such that  $M \sigma_v M'$  (resp.  $N \sigma_v N'$ ), so  $E[V/x] = M[V/x] N[V/x]$  and  $E'[V/x] = M'[V/x] N[V/x]$  (resp.  $E'[V/x] = M[V/x] N'[V/x]$ ); by induction hypothesis,  $M[V/x] \sigma_v M'[V/x]$  (resp.  $N[V/x] \sigma_v N'[V/x]$ ), so  $E[V/x] \sigma_v E'[V/x]$  by the  $@_l$  (resp.  $@_r$ )-rule.

□

**Lemma 139.**

1.  $\beta_v$  (resp.  $\hat{\beta}_v$ ) and  $\sigma_v$  quasi-strongly commute i.e. if  $M \sigma_v N_1$  and  $M \beta_v N_2$  (resp.  $M \hat{\beta}_v N_2$ ) then there exists  $M'$  such that  $N_2 \sigma_v^* M'$  and  $N_1 \beta_v M'$  (resp.  $N_1 \hat{\beta}_v M'$ ).
2.  $\beta_v$  (resp.  $\hat{\beta}_v$ ) and  $\sigma_v$  commute.

PROOF.

1. We prove the statement about  $\beta_v$  by induction on  $M$ . The only interesting cases are:

- if  $M = (\lambda x N)^! V^! L$  with  $M \sigma_1 (\lambda x N L)^! V^! = N_1$  and  $M \beta_v (N[V/x]) L = N_2$ , then  $N_1 \beta_v N_2$  since  $x \notin \text{fv}(L)$ .



- if  $M = (W^!)((\lambda xN)^!V^!)$  with  $M \sigma_3 (\lambda xWN)^!V^! = N_1$  and  $M \beta_v (W^!)N[V/x] = N_2$ , then  $N_1 \beta_v N_2$  since  $x \notin \text{fv}(W)$ .
- if  $M = ((\lambda yP)^!((\lambda xN)^!V^!)L)$  with  $M \sigma_1 (\lambda yPL)^!((\lambda xN)^!V^!) = M_1$  and  $M \beta_v (\lambda yP)^!N[V/x]L = M_2$ , then  $M_1 \beta_v (\lambda yPL)^!N[V/x] = M'$  and  $M_2 \sigma_1 M'$ .
- if  $M = (\lambda xN)^!V^!$  with  $M \sigma_v (\lambda xN)^!V^! = N_1$ ,  $M \beta_v N[V/x] = N_2$  and  $V \sigma_v V'$ , then  $N_1 \beta_v N[V'/x] = M'$  and so  $N_2 \sigma_v^* M'$  by lemma 138.1.
- if  $M = (\lambda xN)^!V^!$  with  $M \sigma_v (\lambda xN')^!V^! = N_1$ ,  $M \beta_v N[V/x] = N_2$  and  $N \sigma_v N'$ , then  $N_1 \beta_v N'[V/x] = M'$  and so  $N_2 \sigma_v M'$  by lemma 138.2.

As regards the statement about  $\hat{\beta}_v$ , it is not proved explicitly because it is enough to observe that in the previous proof whenever the step is  $\hat{\beta}_v$  then we can close the commutation diagram with one  $\hat{\beta}_v$ -reduction step.

2. We prove the following stronger statement, in order to apply the right induction hypothesis: given  $R \in \{\beta_v, \hat{\beta}_v\}$ , if  $L \sigma_v^* N$  and  $L R^m M$  then there exists  $L'$  such that  $M \sigma_v^* L'$  and  $N R^m L'$ . Let  $L \sigma_v^n N$ : the proof is by induction on  $(m, n)$  with the lexicographical order on  $\mathbb{N}^2$ .

If  $m = 0$  or  $n = 0$ , we conclude easily.

Let  $m, n > 0$ : there exist  $N', M'$  such that  $L \sigma N', L R M', N' \sigma_v^{n-1} N$  and  $M' R^{m-1} M$ . By lemma 139.1 applied to  $L$ , there exists  $L''$  such that  $N' R L''$  and  $M \sigma_v^* L''$ . By induction hypothesis applied to  $M'$ , there exists  $M''$  such that  $M \sigma_v^* M''$  and  $L'' R^{m-1} M''$ ; thus  $N' R^m N'$ , so there exists  $L'$  such that  $M'' \sigma_v^* L'$  and  $N R^m L'$  by applying the induction hypothesis to  $N'$ , therefore  $M \sigma_v^* L'$ .  $\square$

**Theorem 140.**  $\beta_{v\sigma}$  and  $\hat{\beta}_{v\sigma}$  are confluent.

PROOF. By proposition 136 and lemmas 137 and 139, since  $\beta_v$  (see [Ehr12]) and  $\hat{\beta}_v$  (see corollary 115.1) are confluent.  $\square$

### 3.3 Simulation of Accattoli and Paolini's calculus and solvability

We present the Accattoli and Paolini's call-by-value  $\lambda$ -calculus with explicit substitutions,  $\lambda_{\text{vsub}}$ , introduced in [AP12]. This calculus can be seen as a merging of two already existing  $\lambda$ -calculi, the Herbelin and Zimmerman's one (a call-by-value  $\lambda$ -calculus with explicit substitutions, see [HZ09]) and the Accattoli and Kesner's one (a call-by-name  $\lambda$ -calculus with explicit substitutions and a very elegant notion of reduction, see [AK10, AK12]).

The following definitions 141, 142 and 145 are exactly the same as in [AP12].

**Definition 141** (Syntax of  $\lambda_{\text{vsub}}$ ). *We define the sets  $\lambda_{\text{vsub}}^{\text{terms}}$  (of  $\lambda_{\text{vsub}}$ -terms) and  $\lambda_{\text{vsub}}^{\text{values}}$  (of  $\lambda_{\text{vsub}}$ -values by mutual induction as follows:*

$$\begin{array}{ll} \lambda_{\text{vsub}}^{\text{term}} & s, t ::= v \mid (s)t \mid s\{t/x\} & \lambda_{\text{vsub}}\text{-terms} \\ \lambda_{\text{vsub}}^{\text{value}} & u, v ::= x \mid \lambda x s & \lambda_{\text{vsub}}\text{-values} \end{array}$$

A constructor of the form  $\{t/x\}$  is an explicit substitution and a term of the form  $s\{t/x\}$  is a term with an explicit substitution. For any  $n \in \mathbb{N}$ , a tuple  $(\{t_1/x_1\}, \dots, \{t_n/x_n\})$  of explicit substitutions is denoted by  $\{t_1/x_1\} \dots \{t_n/x_n\}$ .

Notice that any  $\lambda_{\text{vsub}}$ -value (i.e. a variable or an abstraction) is a  $\lambda_{\text{vsub}}$ -term.

There are two kinds of binder:  $\lambda x t$  and  $t\{u/x\}$ , both binding  $x$  in  $t$ . All  $\lambda_{\text{vsub}}$ -terms are considered up to  $\alpha$ -equivalence. The capture-avoiding substitution of values replacing variables is extended to  $\lambda_{\text{vsub}}$ -terms with explicit substitutions by setting:

$$s\{t/y\}[v/x] = s[v/x]\{t[v/x]/y\}$$

for every  $\lambda_{\text{vsub}}$ -terms  $s$  and  $t$ ,  $\lambda_{\text{vsub}}$ -value  $v$  and variable  $x$  with  $y \notin \text{fv}(v) \cup \{x\}$ .

**Definition 142** ( $\rightarrow_{\lambda_{\text{vsub}}}$ ,  $\rightarrow_{\text{w}}$  and  $\rightarrow_{\text{sw}}$ -reduction). *Let  $R$  be a binary relation on  $\lambda_{\text{vsub}}^{\text{term}}$ .*

*The contextual closure of  $R$  is the binary relation  $R'$  on  $\lambda_{\text{vsub}}^{\text{term}}$  defined by applying, a finite number of times, the following rules:*

$$\begin{array}{c} \frac{s R t}{s R' t} R \quad \frac{s R' s'}{st R' s't} @_l \quad \frac{t R' t'}{st R' st'} @_r \\ \frac{s R' s'}{\lambda x s R' \lambda x s'} \lambda \quad \frac{s R' s'}{s\{t/x\} R' s'\{t/x\}} \text{sub}_l \\ \frac{t R' t'}{s\{t/x\} R' s'\{t/x\}} \text{sub}_r \end{array}$$

*The applicative closure of  $R$  is the binary relation  $R'$  on  $\lambda_{\text{vsub}}^{\text{term}}$  defined by applying, a finite number of times, the following rules:*

$$\begin{array}{c} \frac{s R t}{s R' t} R \quad \frac{s R' s'}{st R' s't} @_l \quad \frac{t R' t'}{st R' st'} @_r \\ \frac{s R' s'}{s\{t/x\} R' s'\{t/x\}} \text{sub}_l \quad \frac{t R' t'}{s\{t/x\} R' s'\{t/x\}} \text{sub}_r \end{array}$$

$\rightarrow_{\lambda_{\text{vsub}}}$  (resp.  $\rightarrow_{\text{w}}$ ), called the  $\rightarrow_{\lambda_{\text{vsub}}}$ -reduction (resp. weak  $\rightarrow_{\lambda_{\text{vsub}}}$ -reduction) is the contextual (resp. applicative) closure of the binary relation  $\mapsto_{\lambda_{\text{vsub}}}$  (resp.  $\mapsto_{\text{w}}$ ) on  $\lambda_{\text{vsub}}^{\text{term}}$  defined by the following rules:

$$\frac{}{(\lambda x s)Lt \mapsto_{\lambda_{\text{vsub}}} s\{t/x\}L} \text{d}\beta \qquad \frac{}{s\{\mathbf{v}L/x\} \mapsto_{\lambda_{\text{vsub}}} s[\mathbf{v}/x]L} \text{sv}$$

where  $L = \{t_1/x_1\} \dots \{t_n/x_n\}$  for some  $n \in \mathbb{N}$  and  $\mathbf{v} \in \lambda_{\text{vsub}}^{\text{value}}$ , moreover  $x_i \notin \text{fv}(t)$  (resp.  $x_i \notin \text{fv}(s)$ ) for every  $1 \leq i \leq n$  in the  $\text{d}\beta$  (resp.  $\text{sv}$ )-rule.

The stratified-weak  $\lambda_{\text{vsub}}$ -reduction is the binary relation  $\rightarrow_{\text{sw}}$  on  $\lambda_{\text{vsub}}^{\text{term}}$  defined by applying, a finite number of times, the following rules:

$$\frac{s \rightarrow_{\text{w}} t}{s \rightarrow_{\text{sw}} t} \text{w} \qquad \frac{s \rightarrow_{\text{sw}} s'}{st \rightarrow_{\text{sw}} s't} \text{@}_1$$

$$\frac{s \rightarrow_{\text{sw}} s'}{\lambda x s \rightarrow_{\text{sw}} \lambda x s'} \lambda \qquad \frac{s \rightarrow_{\text{sw}} s'}{s\{t/x\} \rightarrow_{\text{sw}} s'\{t/x\}} \text{sub}_1$$

The  $\text{d}\beta$ -rule (coming from the call-by-name  $\lambda$ -calculus with explicit substitutions introduced in [AK12]) extend the notion of  $\beta$ -redex: indeed, given some  $\lambda_{\text{vsub}}$ -terms  $s$ ,  $t$  and  $u$ ,  $(\lambda x s)\{t/y\}u \rightarrow_{\lambda_{\text{vsub}}} s\{u/x\}\{t/y\}$  by the  $\text{d}\beta$ -rule. This means that the  $\text{d}\beta$ -rule acts a distance. In the proof-nets representation of  $\lambda_{\text{vsub}}$ -terms this apparent distance is avoided, the  $\text{d}\beta$ -rule is perfectly local from the proof-nets point of view.

The  $\text{sv}$ -rule impose the ‘‘call-by-value’’ constraint in  $\lambda_{\text{vsub}}$ , because only an explicit substitution  $\{\mathbf{v}L/x\}$  (where  $\mathbf{v}$  is a  $\lambda_{\text{vsub}}$ -value and  $L$  is a finite sequence of explicit substitutions) can perform an effective substitution of the occurrences of  $x$  for  $\mathbf{v}$ . The fact that  $s\{\mathbf{v}L/x\} \rightarrow_{\lambda_{\text{vsub}}} s[\mathbf{v}/x]L$  by the  $\text{sv}$ -rule means that also the  $\text{sv}$ -rule acts at a distance.

**Remark 143.** Clearly,  $\rightarrow_{\text{w}} \subseteq \rightarrow_{\text{sw}} \subseteq \rightarrow_{\lambda_{\text{vsub}}}$ .

Stratified-weak  $\lambda_{\text{vsub}}$ -reduction extends weak  $\lambda_{\text{vsub}}$ -reduction allowing reduction under top-level abstractions, which have the important property that cannot be duplicated nor erased.

**Proposition 144.**  $\rightarrow_{\lambda_{\text{vsub}}}$ ,  $\rightarrow_{\text{sw}}$  and  $\rightarrow_{\text{w}}$  are confluent.

PROOF. See corollary 1 and lemma 11 in [AP12].  $\square$

In  $\lambda_{\text{vsub}}$  two terms can have the same behavior and differ only for the position of explicit substitutions, which is not relevant because they do not block  $\rightarrow_{\lambda_{\text{vsub}}}$ -redexes. This is formalized in a precise way by  $\text{o}$ -equivalence on  $\lambda_{\text{vsub}}$ -terms.

**Definition 145** ( $\text{o}$ -equivalence). For every  $i \in \{1, 2, 3, 4\}$ , let  $\sim_{\text{o}_i}$  be the contextual closure of the relation  $\text{o}_i$  defined by the  $\text{o}_i$ -rule:

$$\frac{}{t\{s/x\}\{u/y\} \text{o}_1 t\{u/y\}\{s/x\}} \text{o}_1 \quad \text{where } x \notin \text{fv}(u) \text{ and } y \notin \text{fv}(s)$$

$$\frac{}{tu\{s/x\} \text{o}_2 (tu)\{s/x\}} \text{o}_2 \quad \text{where } x \notin \text{fv}(t)$$

$$\frac{}{t\{s/x\}u \text{o}_3 (tu)\{s/x\}} \text{o}_3 \quad \text{where } x \notin \text{fv}(u)$$

$$\frac{}{t\{s\{u/y\}/x\} \text{o}_4 t\{s/x\}\{u/y\}} \text{o}_4 \quad \text{where } y \notin \text{fv}(t)$$

We set  $\sim_{\circ} = \bigcup_{i=1}^4 \sim_{\circ_i}$ . The  $\circ$ -equivalence is the symmetric and reflexive-transitive closure of  $\sim_{\circ}$ , i.e.  $\equiv_{\circ} = (\sim_{\circ}^T)^*$ .

Remark that  $\equiv_{\circ}$  is an equivalence relation on  $\lambda_{\text{vsub}}$  which allows the commutation of explicit substitutions with every constructor of  $\lambda_{\text{vsub}}$  except abstractions.

We remind a standard notion of rewriting theory and some well-known results about it.

**Definition 146** (Strong bisimulation). *Let  $X$  be a set and let  $\rightarrow_X$  be a binary relation on  $X$ .*

*A strong bisimulation for  $(X, \rightarrow_X)$  is a binary symmetric relation  $\equiv$  on  $X$  such that, for every  $s, s', t \in X$ , if  $s \equiv t$  and  $s \rightarrow_X s'$  then there exists  $t' \in X$  such that  $t \rightarrow_X t'$  and  $s' \equiv t'$ .*

*Given an equivalence relation  $\equiv$  on  $X$ :*

- *we denote by  $\rightarrow_{X/\equiv}$  the binary relation on  $X$  defined by:  $s \rightarrow_{X/\equiv} s'$  iff there exists  $t, u \in X$  such that  $s \equiv t \rightarrow_X u \equiv s'$ ;*
- *we set  $\leftrightarrow_{X/\equiv} = (\rightarrow_X^T \cup \equiv)^*$ ;*
- *$\rightarrow_X$  is Church-Rosser modulo  $\equiv$  if for every  $s, s' \in X$  such that  $s \leftrightarrow_{X/\equiv} s'$ , there exist  $t, t' \in X$  such that  $s \rightarrow_X^* t$ ,  $t \equiv t'$  and  $s' \rightarrow_X^* t'$ .*

**Remark 147.** Let  $X$  be a set, let  $\rightarrow_X$  be a binary relation on  $X$  and let  $\equiv$  be a strong bisimulation for  $(X, \rightarrow_X)$ . If  $\rightarrow_X$  is Church-Rosser modulo  $\equiv$  then  $\rightarrow_X$  is confluent modulo  $\equiv$ , i.e. for every  $s, s', u, u' \in X$  such that  $s \equiv u$ ,  $s \rightarrow_X^* s'$  and  $u \rightarrow_X^* u'$ , there exist  $t, t' \in X$  such that  $s \rightarrow_X^* t$ ,  $t \equiv t'$  and  $s' \rightarrow_X^* t'$ .

**Lemma 148.** *Let  $X$  be a set, let  $\rightarrow_X$  be a binary relation on  $X$  and let  $\equiv$  be a strong bisimulation for  $(X, \rightarrow_X)$  which is an equivalence relation on  $X$ .*

1.  *$\equiv$  can be postponed to  $\rightarrow_X$ , i.e. for every  $s, s' \in X$ , if  $s \rightarrow_{X/\equiv}^* s'$  then there exists  $t \in X$  such that  $s \rightarrow_X^* t \equiv s'$ .*
2. *If  $\rightarrow_X$  is confluent then  $\rightarrow_{X/\equiv}$  is confluent and  $\rightarrow_X$  is Church-Rosser modulo  $\equiv$ .*

PROOF. See for example [Acc11], pp. 86-87. □

Accattoli and Paolini showed that:

**Lemma 149.**  $\equiv_{\circ}$  is a strong bisimulation for both  $(\lambda_{\text{vsub}}^{\text{term}}, \rightarrow_{\lambda_{\text{vsub}}})$  and  $(\lambda_{\text{vsub}}^{\text{term}}, \rightarrow_{\text{sw}})$ .

PROOF. See lemma 12 in [AP12]. □

Therefore, according to proposition 144 and lemmas 148 and 149,  $\equiv_{\circ}$  can be postponed to  $\rightarrow_{\lambda_{\text{vsub}}}$  and  $\rightarrow_{\text{sw}}$ ,  $\rightarrow_{\lambda_{\text{vsub}}/\equiv_{\circ}}$  and  $\rightarrow_{\text{sw}/\equiv_{\circ}}$  are confluent and  $\rightarrow_{\lambda_{\text{vsub}}}$  and  $\rightarrow_{\text{sw}}$  are Church-Rosser modulo  $\equiv_{\circ}$ .

There is a natural way to simulate the  $\lambda_{\text{vsub}}$ -calculus into our  $\Lambda_{\text{CBV}}$ : it is based on the following translation which transforms a  $\lambda_{\text{vsub}}$ -term with an explicit substitution in a  $\beta_{\text{v}}$ -redex in  $\Lambda_{\text{CBV}}$ .

**Definition 150.** *With every  $\lambda_{\text{vsub}}$ -term  $t$  there is associated a term  $(t)^{\diamond} \in \Lambda_{\text{CBV}}$  (also denoted by  $t^{\diamond}$ ) as follows (the definition is by induction on  $t \in \lambda_{\text{vsub}}^{\text{term}}$ ):*

- $(x)^{\diamond} = x^!$ ;
- $(\lambda x t)^{\diamond} = (\lambda x t^{\diamond})^!$ ;
- $(st)^{\diamond} = s^{\diamond}t^{\diamond}$ ;
- $(s\{t/x\})^{\diamond} = (\lambda x s^{\diamond})^!t^{\diamond}$

*With every  $\lambda_{\text{vsub}}$ -value  $\mathbf{v}$  there is associated a value  $(\mathbf{v})^{\blacklozenge} \in \Lambda_{\text{CBV}}$  (also denoted by  $\mathbf{v}^{\blacklozenge}$ ) as follows:*

- $(x)^{\blacklozenge} = x$ ;
- $(\lambda x s)^{\blacklozenge} = \lambda x s^{\blacklozenge}$ .

**Remark 151.** It is immediate to check that:

1. for every  $\mathbf{v} \in \lambda_{\text{vsub}}^{\text{value}}$ , one has  $\mathbf{v}^{\blacklozenge} = (\mathbf{v}^{\blacklozenge})^!$ ;
2. for every  $t \in \lambda_{\text{vsub}}^{\text{term}}$ , one has  $\text{fv}(t) = \text{fv}(t^{\diamond}) = \text{fv}(t^{\blacklozenge})$  (the proof is by straightforward induction on  $t \in \lambda_{\text{vsub}}^{\text{term}}$ ).

**Lemma 152.** *For every  $\lambda_{\text{vsub}}$ -term  $t$ ,  $\lambda_{\text{vsub}}$ -value  $\mathbf{v}$  and variable  $x$ , one has  $(t[\mathbf{v}/x])^{\diamond} = t^{\diamond}[\mathbf{v}^{\blacklozenge}/x]$ .*

PROOF. By induction on the  $\lambda_{\text{vsub}}$ -term  $t$ .

If  $t = x$  then  $(t[\mathbf{v}/x])^{\diamond} = \mathbf{v}^{\blacklozenge} = (\mathbf{v}^{\blacklozenge})^! = t^{\diamond}[\mathbf{v}^{\blacklozenge}/x]$ , by remark 151.

If  $t = y$  for some variable  $y \neq x$ , then  $(t[\mathbf{v}/x])^{\diamond} = y^! = t^{\diamond}[\mathbf{v}^{\blacklozenge}/x]$ .

If  $t = \lambda y s$  for some  $\lambda_{\text{vsub}}$ -term  $s$ , then  $(s[\mathbf{v}/x])^{\diamond} = s^{\diamond}[\mathbf{v}^{\blacklozenge}/x]$  by induction hypothesis, moreover we can suppose without loss of generality that  $y \notin \text{fv}(\mathbf{v}) \cup \{x\}$ , thus  $(t[\mathbf{v}/x])^{\diamond} = (\lambda y (s[\mathbf{v}/x])^{\diamond})^! = (\lambda y s^{\diamond}[\mathbf{v}^{\blacklozenge}/x])^! = t^{\diamond}[\mathbf{v}^{\blacklozenge}/x]$ .

If  $t = s\{u/y\}$  for some  $\lambda_{\text{vsub}}$ -terms  $s$  and  $u$ , then  $(s[\mathbf{v}/x])^{\diamond} = s^{\diamond}[\mathbf{v}^{\blacklozenge}/x]$  and  $(u[\mathbf{v}/x])^{\diamond} = u^{\diamond}[\mathbf{v}^{\blacklozenge}/x]$  by induction hypothesis, moreover we can suppose without loss of generality that  $y \notin \text{fv}(\mathbf{v}) \cup \{x\}$ , hence by lemma 152

$$\begin{aligned} (t[\mathbf{v}/x])^{\diamond} &= (s[\mathbf{v}/x]\{u[\mathbf{v}/x]/y\})^{\diamond} = (\lambda y (s[\mathbf{v}/x])^{\diamond})^!(u[\mathbf{v}/x])^{\diamond} = (\lambda y s^{\diamond}[\mathbf{v}^{\blacklozenge}/x])^!u^{\diamond}[\mathbf{v}^{\blacklozenge}/x] \\ &= ((\lambda y s^{\diamond})^!u^{\diamond})[\mathbf{v}^{\blacklozenge}/x] = t^{\diamond}[\mathbf{v}^{\blacklozenge}/x] \end{aligned}$$

If  $t = su$  for some  $\lambda_{\text{vsub}}$ -terms  $s$  and  $u$ , then  $(s[\mathbf{v}/x])^{\diamond} = s^{\diamond}[\mathbf{v}^{\blacklozenge}/x]$  and  $(u[\mathbf{v}/x])^{\diamond} = u^{\diamond}[\mathbf{v}^{\blacklozenge}/x]$  by induction hypothesis, so  $(t[\mathbf{v}/x])^{\diamond} = s[\mathbf{v}/x]^{\diamond}u[\mathbf{v}/x]^{\diamond} = s^{\diamond}[\mathbf{v}^{\blacklozenge}/x]u^{\diamond}[\mathbf{v}^{\blacklozenge}/x] = t^{\diamond}[\mathbf{v}^{\blacklozenge}/x]$ .  $\square$

In order to simulate in  $\Lambda_{\text{CBV}}$  all the reductions seen in definition 142, we introduce the following notions.

**Definition 153.** The weak  $\beta_{v\sigma}$ -reduction is a binary relation  $\beta_w$  on  $\Lambda_{\text{CBV}}$  defined by applying a finite number of times the following rules:

$$\frac{\frac{M \hat{\beta}_{v\sigma} M'}{M \beta_w N} \hat{\beta}_{v\sigma}}{MN \beta_w M'N} @_l \quad \frac{\frac{M \beta_w M'}{(\lambda x M)! N \beta_w (\lambda x M')! N} \text{red}_l}{N \beta_w N'} @_r$$

The stratified-weak  $\beta_{v\sigma}$ -reduction is a binary relation  $\beta_{\text{sw}}$  on  $\Lambda_{\text{CBV}}$  defined by applying a finite number of times the following rules:

$$\frac{M \beta_w M'}{M \beta_{\text{sw}} N} w \quad \frac{M \beta_{\text{sw}} M'}{MN \beta_{\text{sw}} M'N} @_l \quad \frac{M \beta_{\text{sw}} M'}{\lambda x M \beta_{\text{sw}} \lambda x M'} \lambda$$

$$\frac{V \beta_{\text{sw}} V'}{V! \beta_{\text{sw}} V!} !$$

**Remark 154.** By remark 105, one has  $\beta_w, \beta_{\text{sw}} \subseteq (\Lambda_t \times \Lambda_t) \cup (\Lambda_v \times \Lambda_v)$ . Furthermore,  $\hat{\beta}_{v\sigma} \subseteq \beta_w \subseteq \beta_{\text{sw}} \subseteq \beta_{v\sigma}$ . The converses do not hold.

**Proposition 155** (Simulation of  $\lambda_{\text{vsub}}$  in  $\Lambda_{\text{CBV}}$ ). Let  $s, t \in \lambda_{\text{vsub}}^{\text{term}}$ .

1. If  $s \rightarrow_{\lambda_{\text{vsub}}} t$  then  $s^\diamond \beta_{v\sigma}^* t^\diamond$ .
2. If  $s \rightarrow_w t$  then  $s^\diamond \beta_w^* t^\diamond$ .
3. If  $s \rightarrow_{\text{sw}} t$  then  $s^\diamond \beta_{\text{sw}}^* t^\diamond$ .

PROOF.

1. By induction on the length of the derivation of  $s \rightarrow_{\lambda_{\text{vsub}}} t$ . Let us consider the last rule of this derivation.

If it is the  $d\beta$ -rule, then  $s = (\lambda x u)Lw$  and  $t = u\{w/x\}L$  for some  $\lambda_{\text{vsub}}$ -terms  $u, w$  and tuple of explicit substitutions  $L = \{t_1/x_1\} \dots \{t_n/x_n\}$  with  $n \in \mathbb{N}$ . We can suppose without loss of generality that  $x_i \notin \text{fv}(w)$  for every  $1 \leq i \leq n$ , so

$$s^\diamond = (\lambda x_n \dots (\lambda x_1 (\lambda x u^\diamond)!) t_1^\diamond \dots) t_n^\diamond w^\diamond \sigma_1^n (\lambda x_n \dots (\lambda x_1 (\lambda x u^\diamond)!) w^\diamond) t_1^\diamond \dots) t_n^\diamond = t^\diamond$$

(in particular  $s^\diamond = t^\diamond$  if  $n = 0$ ), therefore  $s^\diamond \hat{\beta}_{v\sigma}^* t^\diamond$  and thus  $s^\diamond \beta_w^* t^\diamond$  (by the  $\hat{\beta}_{v\sigma}$ -rule for  $\beta_w$ ) and  $s^\diamond \beta_{v\sigma}^* t^\diamond$ .

If it is the  $sv$ -rule, then  $s = u\{vL/x\}$  and  $t = u[v/x]L$  for some  $\lambda_{\text{vsub}}$ -term  $u$ ,  $\lambda_{\text{vsub}}$ -value  $v$  and tuple of explicit substitutions  $L = \{t_1/x_1\} \dots \{t_n/x_n\}$  with  $n \in \mathbb{N}$ . We can suppose without loss of generality that  $x_i \notin \text{fv}(u) \cup \{x\}$  for every  $1 \leq i \leq n$ , hence by lemma 152 and remark 151

$$s^\diamond = (\lambda x u^\diamond) ((\lambda x_n \dots (\lambda x_1 v^\diamond)!) t_1^\diamond \dots) t_n^\diamond \sigma_3^n (\lambda x_n \dots (\lambda x_1 (\lambda x u^\diamond)!) (v^\diamond)!) t_1^\diamond \dots) t_n^\diamond$$

$$\beta_v (\lambda x_n \dots (\lambda x_1 u^\diamond [v^\diamond/x]) t_1^\diamond \dots) t_n^\diamond = (\lambda x_n \dots (\lambda x_1 (u\{v/x\})^\diamond) t_1^\diamond \dots) t_n^\diamond = t^\diamond.$$

therefore  $s^\diamond \beta_w^* t^\diamond$  and thus  $s^\diamond \beta_{v\sigma}^* t^\diamond$ .

If it is the  $@_l$ -rule (resp.  $@_r$ -rule) then  $s = uw$  and  $t = u'w$  (resp.  $t = uw'$ ) for some  $\lambda_{v\text{sub}}$ -terms  $u$ ,  $u'$  and  $w$  (resp.  $w'$ ) such that  $u \rightarrow_{\lambda_{v\text{sub}}} u'$  (resp.  $w \rightarrow_{\lambda_{v\text{sub}}} w'$ ). By induction hypothesis,  $u^\diamond \beta_{v\sigma}^* u'^\diamond$  (resp.  $w^\diamond \beta_{v\sigma}^* w'^\diamond$ ), so  $s^\diamond = u^\diamond w^\diamond \beta_{v\sigma}^* u'^\diamond w^\diamond = t^\diamond$  (resp.  $s^\diamond = u^\diamond w^\diamond \beta_{v\sigma}^* u'^\diamond w'^\diamond = t^\diamond$ ) since  $\beta_{v\sigma}^*$  passes to context.

If it is the  $\lambda$ -rule then  $s = \lambda x u$  and  $t = \lambda x u'$  for some  $\lambda_{v\text{sub}}$ -terms  $u$  and  $u'$  such that  $u \rightarrow_{\lambda_{v\text{sub}}} u'$ . By induction hypothesis,  $u^\diamond \beta_{v\sigma}^* u'^\diamond$ , thus  $s^\diamond = (\lambda x u^\diamond)! \beta_{v\sigma}^* (\lambda x u'^\diamond)! = t^\diamond$  since  $\beta_{v\sigma}^*$  passes to context.

If it is the  $sub_l$ -rule (resp.  $sub_r$ -rule) then  $s = u\{w/x\}$  and  $t = u'\{w/x\}$  (resp.  $t = u\{w'/x\}$ ) for some  $\lambda_{v\text{sub}}$ -terms  $u$ ,  $u'$  and  $w$  (resp.  $w'$ ) such that  $u \rightarrow_{\lambda_{v\text{sub}}} u'$  (resp.  $w \rightarrow_{\lambda_{v\text{sub}}} w'$ ). By induction hypothesis,  $u^\diamond \beta_{v\sigma}^* u'^\diamond$  (resp.  $w^\diamond \beta_{v\sigma}^* w'^\diamond$ ), so  $s^\diamond = (\lambda x u^\diamond)! w^\diamond \beta_{v\sigma}^* (\lambda x u'^\diamond)! w^\diamond = t^\diamond$  (resp.  $s^\diamond = (\lambda x u^\diamond)! w^\diamond \beta_{v\sigma}^* (\lambda x u'^\diamond)! w'^\diamond = t^\diamond$ ) since  $\beta_{v\sigma}^*$  passes to context.

2. By induction on the length of the derivation of  $s \rightarrow_w t$ . Let us consider the last rule of this derivation.

If it is the  $d\beta$ - or  $sv$ -rule, then we have seen in the proof of proposition 155.1 that  $s^\diamond \beta_w^* t^\diamond$ .

If it is the  $@_l$ -rule (resp.  $@_r$ -rule) then  $s = uw$  and  $t = u'w$  (resp.  $t = uw'$ ) for some  $\lambda_{v\text{sub}}$ -terms  $u$ ,  $u'$  and  $w$  (resp.  $w'$ ) such that  $u \rightarrow_w u'$  (resp.  $w \rightarrow_w w'$ ). By induction hypothesis,  $u^\diamond \beta_w^* u'^\diamond$  (resp.  $w^\diamond \beta_w^* w'^\diamond$ ), so  $s^\diamond = u^\diamond w^\diamond \beta_w^* u'^\diamond w^\diamond = t^\diamond$  (resp.  $s^\diamond = u^\diamond w^\diamond \beta_w^* u'^\diamond w'^\diamond = t^\diamond$ ) by the  $@_l$ -rule (resp.  $@_r$ -rule) for  $\beta_w$ .

If it is the  $sub_l$ -rule (resp.  $sub_r$ -rule) then  $s = u\{w/x\}$  and  $t = u'\{w/x\}$  (resp.  $t = u\{w'/x\}$ ) for some  $\lambda_{v\text{sub}}$ -terms  $u$ ,  $u'$  and  $w$  (resp.  $w'$ ) such that  $u \rightarrow_w u'$  (resp.  $w \rightarrow_w w'$ ). By induction hypothesis,  $u^\diamond \beta_w^* u'^\diamond$  (resp.  $w^\diamond \beta_w^* w'^\diamond$ ), so  $s^\diamond = (\lambda x u^\diamond)! w^\diamond \beta_w^* (\lambda x u'^\diamond)! w^\diamond = t^\diamond$  (resp.  $s^\diamond = (\lambda x u^\diamond)! w^\diamond \beta_w^* (\lambda x u'^\diamond)! w'^\diamond = t^\diamond$ ) by the  $red_l$ -rule (resp.  $@_r$ -rule) for  $\beta_w$ .

3. By induction on the length of the derivation of  $s \rightarrow_{sw} t$ . Let us consider the last rule of this derivation.

If it is the  $w$ -rule, then  $s^\diamond \beta_w^* t^\diamond$  by 155.2, thus  $s^\diamond \beta_{sw}^* t^\diamond$  by the  $w$ -rule for  $\beta_{sw}$ .

If it is the  $@_l$ -rule then  $s = uw$  and  $t = u'w$  for some  $\lambda_{v\text{sub}}$ -terms  $u$ ,  $u'$  and  $w$  such that  $u \rightarrow_{sw} u'$ . By induction hypothesis,  $u^\diamond \beta_{sw}^* u'^\diamond$ , so  $s^\diamond = u^\diamond w^\diamond \beta_{sw}^* u'^\diamond w^\diamond = t^\diamond$  by the  $@_l$ -rule for  $\beta_{sw}$ .

If it is the  $sub_l$ -rule then  $s = u\{w/x\}$  and  $t = u'\{w/x\}$  for some  $\lambda_{v\text{sub}}$ -terms  $u$ ,  $u'$  and  $w$  such that  $u \rightarrow_{sw} u'$ . By induction hypothesis,  $u^\diamond \beta_{sw}^* u'^\diamond$ , so  $s^\diamond = (\lambda x u^\diamond)! w^\diamond \beta_{sw}^* (\lambda x u'^\diamond)! w^\diamond = t^\diamond$  by the  $\lambda$ -rule,  $!$ -rule and  $@_l$ -rule for  $\beta_{sw}$ .

If it is the  $\lambda$ -rule then  $s = \lambda x u$  and  $t = \lambda x u'$  for some  $\lambda_{\text{vsub}}$ -terms  $u$  and  $u'$  such that  $u \rightarrow_{\text{sw}} u'$ . By induction hypothesis,  $u^\diamond \beta_{\text{sw}}^* u'^\diamond$ , so  $s^\diamond = \lambda x u^\diamond \beta_{\text{sw}}^* \lambda x u'^\diamond = t^\diamond$  by the  $\lambda$ -rule for  $\beta_{\text{sw}}$ .  $\square$

**Remark 156.**  $s \rightarrow_w t$  does not implies that  $s^\diamond \hat{\beta}_{\text{v}\sigma}^* t^\diamond$ . For instance, take  $u = (z_1)z_2 \in \lambda_{\text{vsub}}^{\text{term}}$  and  $s = x_1\{y_1\{u/y_2\}/x_2\} \in \lambda_{\text{vsub}}^{\text{term}}$  where  $x_1, x_2, y_1, y_2, z_1, z_2$  are pairwise distinct variables: then  $u^\diamond = (z_1^!)z_2^!$ ,  $s \rightarrow_w x_1\{u/y_2\} = t$  and

$$s^\diamond = (\lambda x_2 x_1^!)((\lambda y_2 y_1^!)u^\diamond) \sigma_3 (\lambda y_2 (\lambda x_2 x_1^!)y_1^!)u^\diamond = M$$

where  $M$  is  $\hat{\beta}_{\text{v}\sigma}$ -normal but  $t^\diamond = (\lambda y_2 x_1^!)u^\diamond \neq M$ . On the contrary,  $M \beta_w (\lambda y_2 x_1^!)u^\diamond = t^\diamond$  thanks to  $\text{red}_1$ -rule.

By means of “forgetful functor”  $(\ )^{\text{F}}$  (see p. 65) and  $\text{o}$ -equivalence,  $\rightarrow_{\lambda_{\text{vsub}}}$ -reduction can simulate the  $\beta_{\text{v}\sigma}$ -reduction.

**Remark 157.** For every  $E \in \Lambda_{\text{CBV}}$  one has  $\text{fv}(E) = \text{fv}(E^{\text{F}})$  (the proof is by straightforward induction on  $E \in \Lambda_{\text{CBV}}$ ).

**Lemma 158.** For every term  $M$ , value  $V$  and variable  $x$ , one has  $(M[V/x])^{\text{F}} = M^{\text{F}}[V^{\text{F}}/x]$ .

PROOF. By induction on  $M \in \Lambda_{\text{t}}$ .

If  $M = x$  then  $M^{\text{F}} = x$ , thus  $(M[V/x])^{\text{F}} = V^{\text{F}} = M^{\text{F}}[V^{\text{F}}/x]$ .

If  $M = y$  for some variable  $x \neq y$ ,  $M^{\text{F}} = y$ , thus  $(M[V/x])^{\text{F}} = y = M^{\text{F}}[V^{\text{F}}/x]$ .

If  $M = \lambda y N$  for some term  $N$ , then we can suppose without loss of generality that  $y \notin \text{fv}(V) \cup \{x\}$ ; by induction hypothesis,  $(N[V/x])^{\text{F}} = N^{\text{F}}[V^{\text{F}}/x]$ , thus  $(M[V/x])^{\text{F}} = (\lambda y N[V/x])^{\text{F}} = \lambda y (N[V/x])^{\text{F}} = \lambda y N^{\text{F}}[V^{\text{F}}/x] = M^{\text{F}}[V^{\text{F}}/x]$ .

If  $M = W^!$  for some value  $W$ , then  $(W[V/x])^{\text{F}} = W^{\text{F}}[V^{\text{F}}/x]$  by induction hypothesis, so  $(M[V/x])^{\text{F}} = ((W[V/x])^!)^{\text{F}} = (W[V/x])^{\text{F}} = W^{\text{F}}[V^{\text{F}}/x] = M^{\text{F}}[V^{\text{F}}/x]$ .

If  $M = NL$  for some terms  $N$  and  $L$ , then  $(N[V/x])^{\text{F}} = N^{\text{F}}[V^{\text{F}}/x]$  and  $(L[V/x])^{\text{F}} = L^{\text{F}}[V^{\text{F}}/x]$ , hence  $(M[V/x])^{\text{F}} = (N[V/x]L[V/x])^{\text{F}} = (N[V/x])^{\text{F}}(L[V/x])^{\text{F}} = N^{\text{F}}[V^{\text{F}}/x]L^{\text{F}}[V^{\text{F}}/x] = M^{\text{F}}[V^{\text{F}}/x]$ .  $\square$

**Lemma 159.** Let  $E, E' \in \Lambda_{\text{CBV}}$ .

1. If  $E \beta_v E'$  then  $E^{\text{F}} \rightarrow_{\lambda_{\text{vsub}}}^+ E'^{\text{F}}$ .
2. If  $E \sigma_v E'$  then  $E^{\text{F}} \leftrightarrow_{\lambda_{\text{vsub}}/\equiv_{\text{o}}} E'^{\text{F}}$ .
3. If  $E \beta_{\text{v}\sigma}^* E'$  then  $E^{\text{F}} \leftrightarrow_{\lambda_{\text{vsub}}/\equiv_{\text{o}}} E'^{\text{F}}$ .
4. For every  $M \in \Lambda_{\text{CBV}}$ , if  $M \beta_{\text{v}\sigma}^* (\lambda x x^!)$  then  $M^{\text{F}} \rightarrow_{\lambda_{\text{vsub}}}^* \lambda x x$ .

PROOF.



1. By induction on the derivation of  $E \beta_{\mathbf{v}} E'$ . Let us consider the last rule of this derivation.

If it is the  $\beta$ -rule, then  $E = (\lambda x N)^! V^!$  and  $E' = N[V/x]$ , hence  $E^{\mathbf{F}} = (\lambda x N^{\mathbf{F}}) V^{\mathbf{F}} \rightarrow_{\lambda_{\mathbf{vsub}}} N^{\mathbf{F}}\{V^{\mathbf{F}}/x\} \rightarrow_{\lambda_{\mathbf{vsub}}} N^{\mathbf{F}}[V^{\mathbf{F}}/x] = (N[V/x])^{\mathbf{F}} = E'^{\mathbf{F}}$  by  $\mathbf{d}\beta$ - and  $\mathbf{sv}$ -rule and lemma 158.

If it is the  $\textcircled{\_}$ (resp.  $\textcircled{\_}_r$ )-rule, then  $E = NL$  and  $E' = N'L$  (resp.  $E' = NL'$ ) for some terms  $N$ ,  $L$  and  $N'$  (resp.  $L'$ ) such that  $N \beta_{\mathbf{v}} N'$  (resp.  $L \beta_{\mathbf{v}} L'$ ). By induction hypothesis  $N^{\mathbf{F}} \rightarrow_{\lambda_{\mathbf{vsub}}}^+ N'^{\mathbf{F}}$  (resp.  $L^{\mathbf{F}} \rightarrow_{\lambda_{\mathbf{vsub}}}^+ L'^{\mathbf{F}}$ ), so  $E^{\mathbf{F}} = N^{\mathbf{F}} L^{\mathbf{F}} \rightarrow_{\lambda_{\mathbf{vsub}}}^+ N'^{\mathbf{F}} L^{\mathbf{F}} = E'^{\mathbf{F}}$  (resp.  $E^{\mathbf{F}} = N^{\mathbf{F}} L^{\mathbf{F}} \rightarrow_{\lambda_{\mathbf{vsub}}}^+ N^{\mathbf{F}} L'^{\mathbf{F}} = E'^{\mathbf{F}}$ ).

If it is the  $\lambda$ -rule, then  $E = \lambda x N$  and  $E' = \lambda x N'$  for some terms  $N$  and  $N'$  such that  $N \beta_{\mathbf{v}} N'$ . By induction hypothesis  $N^{\mathbf{F}} \rightarrow_{\lambda_{\mathbf{vsub}}}^+ N'^{\mathbf{F}}$ , thus  $E^{\mathbf{F}} = \lambda x N^{\mathbf{F}} \rightarrow_{\lambda_{\mathbf{vsub}}}^+ \lambda x N'^{\mathbf{F}} = E'^{\mathbf{F}}$ .

If it is the  $!$ -rule then  $E = V^!$  and  $E' = V'^!$  for some values  $V$  and  $V'$  such that  $V \beta_{\mathbf{v}} V'$ . By induction hypothesis  $V^{\mathbf{F}} \rightarrow_{\lambda_{\mathbf{vsub}}}^+ V'^{\mathbf{F}}$ , hence  $E^{\mathbf{F}} = V^{\mathbf{F}} \rightarrow_{\lambda_{\mathbf{vsub}}}^+ V'^{\mathbf{F}} = E'^{\mathbf{F}}$ .

2. By induction on the derivation of  $E \sigma_{\mathbf{v}} E'$ . Let us consider the last rule of this derivation.

If it is the  $\sigma_1$ -rule, then  $E = (\lambda x M)^N L$  and  $E' = (\lambda x M L)^! N$  for some terms  $M$ ,  $N$  and  $L$  with  $x \notin \text{fv}(L) = \text{fv}(L^{\mathbf{F}})$ , thus  $E^{\mathbf{F}} = (\lambda x M^{\mathbf{F}}) N^{\mathbf{F}} L^{\mathbf{F}} \rightarrow_{\lambda_{\mathbf{vsub}}} M^{\mathbf{F}}\{N^{\mathbf{F}}/x\} L^{\mathbf{F}} \equiv_{\circ} (M^{\mathbf{F}} L^{\mathbf{F}})\{N^{\mathbf{F}}/x\} \lambda_{\mathbf{vsub}} \leftarrow (\lambda x M^{\mathbf{F}} L^{\mathbf{F}}) N^{\mathbf{F}} = E'^{\mathbf{F}}$ . Therefore  $E^{\mathbf{F}} \leftrightarrow_{\lambda_{\mathbf{vsub}}/\equiv_{\circ}} E'^{\mathbf{F}}$ .

If it is the  $\sigma_3$ -rule, then  $E = (V^!)((\lambda x L)^!) N$  and  $E' = (\lambda x V^! L)^! N$  for some terms  $N$  and  $L$  and value  $V$  with  $x \notin \text{fv}(V) = \text{fv}(V^{\mathbf{F}})$ , hence  $E^{\mathbf{F}} = (V^{\mathbf{F}})(\lambda x L^{\mathbf{F}}) N^{\mathbf{F}} \rightarrow_{\lambda_{\mathbf{vsub}}} V^{\mathbf{F}} L^{\mathbf{F}}\{N^{\mathbf{F}}/x\} \equiv_{\circ} (V^{\mathbf{F}} L^{\mathbf{F}})\{N^{\mathbf{F}}/x\} \lambda_{\mathbf{vsub}} \leftarrow (\lambda x V^{\mathbf{F}} L^{\mathbf{F}}) N^{\mathbf{F}} = E'^{\mathbf{F}}$ . Therefore  $E^{\mathbf{F}} \leftrightarrow_{\lambda_{\mathbf{vsub}}/\equiv_{\circ}} E'^{\mathbf{F}}$ .

If it is the  $\textcircled{\_}$ (resp.  $\textcircled{\_}_r$ )-rule, then  $E = NL$  and  $E' = N'L$  (resp.  $E' = NL'$ ) for some terms  $N$ ,  $L$  and  $N'$  (resp.  $L'$ ) such that  $N \sigma_{\mathbf{v}} N'$  (resp.  $L \sigma_{\mathbf{v}} L'$ ). By induction hypothesis  $N^{\mathbf{F}} \leftrightarrow_{\lambda_{\mathbf{vsub}}/\equiv_{\circ}} N'^{\mathbf{F}}$  (resp.  $L^{\mathbf{F}} \leftrightarrow_{\lambda_{\mathbf{vsub}}/\equiv_{\circ}} L'^{\mathbf{F}}$ ), so  $M^{\mathbf{F}} = N^{\mathbf{F}} L^{\mathbf{F}} \leftrightarrow_{\lambda_{\mathbf{vsub}}/\equiv_{\circ}} N'^{\mathbf{F}} L^{\mathbf{F}} = M'^{\mathbf{F}}$  (resp.  $E^{\mathbf{F}} = N^{\mathbf{F}} L^{\mathbf{F}} \leftrightarrow_{\lambda_{\mathbf{vsub}}/\equiv_{\circ}} N^{\mathbf{F}} L'^{\mathbf{F}} = E'^{\mathbf{F}}$ ).

If it is the  $\lambda$ -rule, then  $E = \lambda x N$  and  $E' = \lambda x N'$  for some terms  $N$  and  $N'$  such that  $N \sigma_{\mathbf{v}} N'$ . By induction hypothesis  $N^{\mathbf{F}} \leftrightarrow_{\lambda_{\mathbf{vsub}}/\equiv_{\circ}} N'^{\mathbf{F}}$ , thus  $E^{\mathbf{F}} = \lambda x N^{\mathbf{F}} \leftrightarrow_{\lambda_{\mathbf{vsub}}/\equiv_{\circ}} \lambda x N'^{\mathbf{F}} = E'^{\mathbf{F}}$ .

If it is the  $!$ -rule then  $E = V^!$  and  $E' = V'^!$  for some values  $V$  and  $V'$  such that  $V \beta_{\mathbf{v}} V'$ . By induction hypothesis  $V^{\mathbf{F}} \leftrightarrow_{\lambda_{\mathbf{vsub}}/\equiv_{\circ}} V'^{\mathbf{F}}$ , hence  $E^{\mathbf{F}} = V^{\mathbf{F}} \leftrightarrow_{\lambda_{\mathbf{vsub}}/\equiv_{\circ}} V'^{\mathbf{F}} = E'^{\mathbf{F}}$ .

3. By induction on the number  $n \in \mathbb{N}$  of steps of the  $\beta_{v\sigma}$ -reduction from  $E$  to  $E'$ . If  $n = 0$  then  $E = E'$  and so  $E^F = E'^F$ , thus  $E^F \leftrightarrow_{\lambda_{v\text{sub}}/\equiv_{\circ}} E'^F$ . If  $n > 0$  then there exists  $E'' \in \Lambda_{\text{CBV}}$  such that  $E \beta_{v\sigma}^n E'' \beta_{v\sigma} E'$ ; by induction hypothesis  $E^F \leftrightarrow_{\lambda_{v\text{sub}}/\equiv_{\circ}} E''^F$ ; if  $E'' \beta_v E'$  then  $E'^F \rightarrow_{\lambda_{v\text{sub}}}^+ E''^F$  by lemma 159.1, otherwise  $E''^F \leftrightarrow_{\lambda_{v\text{sub}}/\equiv_{\circ}} E'^F$  by lemma 159.2; in any case  $E^F \leftrightarrow_{\lambda_{v\text{sub}}/\equiv_{\circ}} E'^F$ .
4. By lemma 159.3  $M^F \leftrightarrow_{\lambda_{v\text{sub}}/\equiv_{\circ}} (\lambda x x^!)^!F = \lambda x x$ . By lemmas 148.2 and 149, there exist  $t, t' \in \lambda_{v\text{sub}}^{\text{term}}$  such that  $M^F \rightarrow_{\lambda_{v\text{sub}}}^* t$ ,  $t \equiv_{\circ} t'$  and  $\lambda x x \rightarrow_{\lambda_{v\text{sub}}}^* t'$ . As  $\lambda x x$  is  $\rightarrow_{\lambda_{v\text{sub}}}$ -normal, one has  $t' = \lambda x x$ , thus  $t \equiv_{\circ} \lambda x x$  that implies  $t = \lambda x x$ . Therefore  $M^F \rightarrow_{\lambda_{v\text{sub}}}^* \lambda x x$ .  $\square$

### 3.4 From terms to trees

It is more natural to study  $\hat{\beta}_{v!}$ -,  $\hat{\beta}_{vr}$ -,  $\hat{\beta}_v$ - and  $\hat{\beta}_{vt}$ -reductions and the decompositions seen in remark 100.2, theorem 118, corollaries 119.2 and 119.5 and proposition 121 by means of labeled full binary trees. This analysis, perhaps implicit in some publications on call-by-value  $\lambda$ -calculus, has never been made explicitly and reveals deep symmetries of this calculus which can be seen from a more “geometrical” point of view.

#### 3.4.1 Syntax of applicative trees

**Definition 160** (Quasi-leaf, applicative tree). *We denote by  $\mathcal{T}_{\textcircled{a}}$  the set of full binary trees whose internal nodes are labeled by  $\textcircled{a}$  and whose leaves are labeled by terms of the shape  $V^!$  with  $V \in \Lambda_v$ .*

*A quasi-leaf is an element of  $\mathcal{T}_{\textcircled{a}}$  whose root is a  $\textcircled{a}$ -node having two leaves of  $T$  as children. For every term  $T \in \mathcal{T}_{\textcircled{a}}$ , a quasi-leaf of  $T$  is a subtree of  $T$  which is a quasi-leaf.*

*With every term  $M$  is associated  $\text{app}(M) \in \mathcal{T}_{\textcircled{a}}$ , called applicative tree of  $M$ , defined by induction on  $M$  as follows:*

- $\text{app}(V^!)$  consists of a leaf labeled by  $V^!$ ;
- $\text{app}(MN)$  consists of a node labeled by  $\textcircled{a}$  whose left (resp. right) child is  $\text{app}(M)$  (resp.  $\text{app}(N)$ ).

*With every  $T \in \mathcal{T}_{\textcircled{a}}$  is associated  $\text{term}(T) \in \Lambda_v$ , called term of  $T$ , defined by induction on  $T$  as follows:*

- if  $T$  consist simply of a leaf labeled by  $V^!$  for some value  $V$  then  $\text{term}(T) = V^!$ ;
- if  $T$  consists of a node labeled by  $\textcircled{a}$  whose left (resp. right) child is  $\text{term}(T_1)$  (resp.  $\text{term}(T_2)$ ), then  $\text{term}(T) = (\text{term}(T_1))\text{term}(T_2)$ .

**Remark 161.**

1. Every leaf of every  $T \in \mathcal{T}_{\textcircled{a}}$  is labeled by a term of the shape  $V^!$  for some value  $V$ .
2. For every terms  $M$  and  $N$ ,  $\mathbf{app}(MN)$  has at least one quasi-leaf: this is a well-known result on full binary trees having more than one node. As a consequence, for every term  $M$ ,  $\mathbf{app}(M)$  is either a leaf and in this case  $M = V^!$  for some value  $V$ , or such that its root has two children and each sub-trees of  $\mathbf{app}(M)$  contains a quasi-leaf.
3. It is immediate to check that the two functions  $\mathbf{app}$  and  $\mathbf{term}$  are inverses of each other: for every  $T \in \mathcal{T}_{\textcircled{a}}$  and term  $M$ ,  $\mathbf{app}(\mathbf{term}(T)) = T$  and  $\mathbf{term}(\mathbf{app}(M)) = M$ . So each element of  $\mathcal{T}_{\textcircled{a}}$  is the applicative tree of some term and each term is uniquely determined by its applicative tree.

**Definition 162.** Let  $T \in \mathcal{T}_{\textcircled{a}}$ . With every node  $n$  of  $T$  is associated a finite sequence  $\mathbf{pos}_T(n)$  of elements of  $\{l, r\}$  as follows (the definition is by induction on  $T$ ):

- if the root of  $T$  is a leaf, then  $\mathbf{pos}_T(n) = \emptyset$ ;
- if the root of  $T$  is not a leaf and if  $T_l$  (resp.  $T_r$ ) is the left (resp. right) child of  $n$ , then for every node  $m$  in  $T_l$  (resp.  $T_r$ )  $\mathbf{pos}_T(m) = l \cdot \mathbf{pos}_{T_l}(m)$  (resp.  $\mathbf{pos}_T(m) = r \cdot \mathbf{pos}_{T_r}(m)$ ).

For every subtree  $T'$  of  $T$  we set  $\mathbf{pos}_T(T') = \mathbf{pos}_T(n)$  where  $n$  is the root of  $T'$ .

**Remark 163.** Given  $T \in \mathcal{T}_{\textcircled{a}}$ ,  $\mathbf{pos}_T$  is an injection: if  $\mathbf{pos}_T(n) = \mathbf{pos}_T(m)$  (resp.  $\mathbf{pos}_T(T_1) = \mathbf{pos}_T(T_2)$ ) then  $n = m$  (resp.  $T_1 = T_2$ ), for any nodes  $m, n$  (resp. subtrees  $T_1, T_2$ ) of  $T$ ; this is a consequence of acyclicity of trees;

Given a term  $M$ , we can localize uniquely all its subterms thanks to  $\mathbf{pos}_{\mathbf{app}(M)}^{-1}$ .

WHY?

We recall a well-known results on trees.

**Proposition 164.** Let  $T$  be a tree.

1. Given two subterms  $T_1$  and  $T_2$  of  $T$ , either  $T_1$  is a subtree of  $T_2$ , or  $T_2$  is a subtree of  $T_1$ , or  $T_1$  and  $T_2$  are disjoint.
2. Let  $<_T$  be the binary relation on the subtrees of  $T$  defined by:  $R <_T S$  ( $R$  is on the left of  $S$  in  $T$ ) iff  $R$  and  $S$  are subterms of  $T$  such that  $\mathbf{pos}_T(R) = (r_1, \dots, r_{n_R})$ ,  $\mathbf{pos}_T(S) = (s_1, \dots, s_{n_S})$  with  $n_R, n_S \in \mathbb{N}$ ,  $r_i, s_j \in \{l, r\}$  for any  $1 \leq i \leq n_R$  and  $1 \leq j \leq n_S$ , and there exists  $m \leq n_R, n_S$  such that  $r_m = l$ ,  $s_m = r$  and  $r_i = s_i$  for every  $1 \leq i < m$ . Then  $<_T$  is an order relation on the disjoint subtrees of  $T$ .

PROOF.  $T$  is a acyclic and connected graph, so we can conclude.  $\square$

Proposition 164.2 says that a natural order relation “from left to right” is definable on the disjoint subtrees of  $T$ .

**Remark 165.** Proposition 164 has as one of its consequences that for every  $T \in \mathcal{T}_{\text{@}}$ , if  $T_1$  and  $T_2$  are two distinct quasi-leaves of  $T$  then  $T_1$  are  $T_2$  are disjoint and either  $T_1$  is on the left of  $T_2$  or  $T_2$  is on the left of  $T_1$ . This order relation is on the quasi-leaves of  $T$  is well-founded.

This notion of order is useful for the following definition.

**Definition 166** ( $\hat{\beta}_v$ -redex on  $\mathcal{T}_{\text{@}}$ ). A  $\hat{\beta}_v$ -redex on  $\mathcal{T}_{\text{@}}$  is an element of  $\mathcal{T}_{\text{@}}$  which is a quasi-leaf whose left child is labeled by  $(\lambda x M)^!$  for some  $M \in \Lambda_t$ . Let  $T$  be a  $\hat{\beta}_v$ -redex on  $\mathcal{T}_{\text{@}}$  whose left (resp. right) child is labeled by  $(\lambda x M)^!$  (resp.  $V^!$ ) for some term  $M$  (resp. value  $V$ ): the contractum of  $T$  is  $\text{app}(M[V/x])$ .

Let  $T \in \mathcal{T}_{\text{@}}$ . A  $\hat{\beta}_v$ -redex in  $T$  is an occurrence of subtree of  $T$  which is a  $\hat{\beta}_v$ -redex on  $\mathcal{T}_{\text{@}}$ . If  $T'$  is both the leftmost (resp. rightmost) quasi-leaf of  $T$  and a  $\hat{\beta}_v$ -redex on  $\mathcal{T}_{\text{@}}$ ,  $T'$  is the  $\hat{\beta}_{vl}$ - (resp.  $\hat{\beta}_{vr}$ -)redex of  $T$ .

We say that  $T$  contains a  $\hat{\beta}_v$ - (resp.  $\hat{\beta}_{vl}$ -;  $\hat{\beta}_{vr}$ -)redex if there is a  $\hat{\beta}_v$ - (resp.  $\hat{\beta}_{vl}$ -;  $\hat{\beta}_{vr}$ -)redex in  $T$ .

Some elements of  $\mathcal{T}_{\text{@}}$  might contains a  $\hat{\beta}_v$ -redex without having neither the  $\hat{\beta}_{vl}$ -redex nor the  $\hat{\beta}_{vr}$ -redex, for example  $\text{app}(\lambda x_1 x_2. (M)z^!)(y_1^!y_2^!)$  where  $M = (\lambda x N)^!$  for some term  $N$ .

Formally, the notions of  $\hat{\beta}_v$ -,  $\hat{\beta}_{vl}$ - and  $\hat{\beta}_{vr}$ -redex on  $\mathcal{T}_{\text{@}}$  are distinct from the notions of  $\hat{\beta}_v$ -,  $\hat{\beta}_{vl}$ - and  $\hat{\beta}_{vr}$ -redex on terms, but actually they are strictly correlated. For example, if  $M = ((\lambda x_1 x_1^!)^!y^!)((\lambda x_2 x_2^!)^!)z^!$ ,  $\text{app}(M)$  contains exactly two  $\hat{\beta}_v$ -redexes,  $\text{app}((\lambda x_1 x_1^!)^!y^!)$  (the  $\hat{\beta}_{vl}$ -redex of  $\text{app}(M)$ ) and  $\text{app}((\lambda x_2 x_2^!)^!z^!)$  (the  $\hat{\beta}_{vr}$ -redex of  $\text{app}(M)$ ); but  $M$  also contains exactly two  $\hat{\beta}_v$ -redexes,  $(\lambda x_1 x_1^!)^!y^!$  (the  $\hat{\beta}_{vl}$ -redex of  $M$ ) and  $(\lambda x_2 x_2^!)^!z^!$  (the  $\hat{\beta}_{vr}$ -redex of  $M$ ).

**Lemma 167.** Let  $M, N$  be terms:  $\text{app}(N)$  is a  $\hat{\beta}_v$ -redex in  $\text{app}(M)$  iff  $N$  is a  $\hat{\beta}_v$ -redex in  $M$ .

PROOF. By induction on the term  $M$ .

If  $M = V^!$  for some value  $V$ , then there is no  $\hat{\beta}_v$ -redex in  $M$  and  $\text{app}(M)$  consists only of a leaf labeled by  $V^!$ , so there is no  $\hat{\beta}_{vt}$ -redex in  $\text{app}(M)$ .

If  $M = M_1 M_2$  for some terms  $M_1, M_2$ , then there are three cases.

- $N$  (resp.  $\text{app}(N)$ ) is a  $\hat{\beta}_v$ - (resp.  $\hat{\beta}_{vt}$ -)redex in  $M_1$  (resp.  $\text{app}(M_1)$ ): by induction hypothesis,  $\text{app}(N)$  (resp.  $N$ ) is a  $\hat{\beta}_v$ -redex in  $\text{app}(M_1)$  (resp.  $M_1$ ); as  $\text{app}(M)$  consists of a node whose left child is  $\text{app}(M_1)$  (resp. as  $M = M_1 M_2$ ), then  $\text{app}(N)$  (resp.  $N$ ) is a  $\hat{\beta}_v$ -redex of  $\text{app}(M)$  (resp.  $M$ ).

- $N$  (resp.  $\mathbf{app}(N)$ ) is a  $\hat{\beta}_v$ - (resp.  $\hat{\beta}_{vt}$ -)redex in  $M_2$  (resp.  $\mathbf{app}(M_2)$ ): by induction hypothesis,  $\mathbf{app}(N)$  (resp.  $N$ ) is a  $\hat{\beta}_v$ -redex in  $\mathbf{app}(M_2)$  (resp.  $M_2$ ); as  $\mathbf{app}(M)$  consists of a node whose left child is  $\mathbf{app}(M_2)$  (resp. as  $M = M_1M_2$ ), then  $\mathbf{app}(N)$  (resp.  $N$ ) is a  $\hat{\beta}_v$ -redex of  $\mathbf{app}(M)$  (resp.  $M$ ).
- $M = N = (\lambda xN')^!V^!$  for some term  $N'$  and value  $V$ :  $\mathbf{app}(M) = \mathbf{app}(N)$  consists of a quasi-leaf whose left leaf is labeled by  $(\lambda xN')^!$  and whose right leaf is labeled by  $V^!$ , hence  $\mathbf{app}(M)$  is a  $\hat{\beta}_{vt}$ -redex in  $\mathbf{app}(M)$ .  $\square$

### 3.4.2 Some reductions on applicative trees

**Definition 168** ( $\hat{\beta}_v$ -,  $\hat{\beta}_{vl}$ -,  $\hat{\beta}_{vr}$ - and  $\hat{\beta}_{vt}$ -reduction on  $\mathcal{T}_{\textcircled{a}}$ ). *We define a relation  $\hat{\beta}_v \subseteq \mathcal{T}_{\textcircled{a}} \times \mathcal{T}_{\textcircled{a}}$ , called  $\hat{\beta}_v$ -reduction or weak  $\beta_v$ -reduction on  $\mathcal{T}_{\textcircled{a}}$ , as follows:  $T \hat{\beta}_v T'$  if  $T, T' \in \mathcal{T}_{\textcircled{a}}$  and  $T'$  is obtained from  $T$  by replacing a  $\hat{\beta}_v$ -redex in  $T$  with its contractum.*

*We define a relation  $\hat{\beta}_{vl}$  (resp.  $\hat{\beta}_{vr}$ )  $\subseteq \mathcal{T}_{\textcircled{a}} \times \mathcal{T}_{\textcircled{a}}$ , called  $\hat{\beta}_{vl}$  (resp.  $\hat{\beta}_{vr}$ )-reduction or leftmost (resp. rightmost) weak  $\beta_v$ -reduction on  $\mathcal{T}_{\textcircled{a}}$ , as follows:  $T \hat{\beta}_{vl} \hat{\beta}_{vr} T'$  if  $T, T' \in \mathcal{T}_{\textcircled{a}}$  and  $T'$  is obtained from  $T$  by replacing the  $\hat{\beta}_{vl}$  (resp.  $\hat{\beta}_{vr}$ )-redex in  $T$  (if any) with its contractum.*

*We define a relation  $\hat{\beta}_{vt} \subseteq \mathcal{T}_{\textcircled{a}} \times \mathcal{T}_{\textcircled{a}}$ , called  $\hat{\beta}_{vt}$ -reduction or turbo weak  $\beta_v$ -reduction on  $\mathcal{T}_{\textcircled{a}}$ , as follows:  $T \hat{\beta}_{vt} T'$  if  $T, T' \in \mathcal{T}_{\textcircled{a}}$  and  $T'$  is obtained from  $T$  by replacing each  $\hat{\beta}_v$ -redex in  $T$  with its contractum.*

**Theorem 169.** *Let  $T, T' \in \mathcal{T}_{\textcircled{a}}$ .*

- *If  $T \hat{\beta}_v T'$  then  $\mathbf{term}(T) \hat{\beta}_v \mathbf{term}(T')$ .*
- *If  $T \hat{\beta}_{vl}$  (resp.  $\hat{\beta}_{vr}$ )  $T'$  then  $\mathbf{term}(T) \hat{\beta}_{vl}$  (resp.  $\hat{\beta}_{vr}$ )  $\mathbf{term}(T')$ .*
- *If  $T \hat{\beta}_{vt} T'$  then  $\mathbf{term}(T) \hat{\beta}_{vt} \mathbf{term}(T')$ .*

## 3.5 Value Böhm trees

We extend the call-by-value  $\lambda$ -calculus  $\Lambda_{\text{CBV}}$  by adding a constant  $\Omega$ .

**Definition 170.** *We define the elements of the sets  $\Lambda_t^\Omega$  ( $\Omega$ -terms),  $\Lambda_v^\Omega$  ( $\Omega$ -values) and  $\Lambda_{\text{CBV}}^\Omega$  ( $\Omega$ -expressions) by mutual induction as follows:*

$$\begin{array}{lll}
 \Lambda_t^\Omega & L, M, N ::= \Omega \mid (M)N \mid (V)^! & \Omega\text{-terms} \\
 \Lambda_v^\Omega & U, V, W ::= x \mid \lambda xM & \Omega\text{-values} \\
 \Lambda_{\text{CBV}}^\Omega & D, E, F ::= M \mid V & \Omega\text{-expressions}
 \end{array}$$

The constant  $\Omega$  has to be considered as a closed term, in particular  $\Omega[V/x] := \Omega$  for every value  $V$  and variable  $x$ .

**Definition 171.** We define a relation  $\beta_{v\Omega} \subseteq (\Lambda_t^\Omega \times \Lambda_t^\Omega) \cup (\Lambda_v^\Omega \times \Lambda_v^\Omega)$ , called  $\beta_{v\Omega}$ -reduction, by the following rules:

$$\frac{}{(\lambda x M)!V! \beta_{v\Omega} M[V/x]} \beta \quad \frac{M \beta_{v\Omega} M'}{MN \beta_{v\Omega} M'N} @_l$$

$$\frac{N \beta_{v\Omega} N'}{MN \beta_{v\Omega} MN'} @_r$$

$$\frac{M \beta_{v\Omega} M'}{\lambda x M \beta_{v\Omega} \lambda x M'} \lambda \quad \frac{V \beta_{v\Omega} W}{V! \beta_{v\Omega} W!} ! \quad \frac{}{(\Omega)M \beta_{v\Omega} \Omega} \Omega_l$$

$$\frac{}{(M)\Omega \beta_{v\Omega} \Omega} \Omega_r$$

We prove now a confluence property for this calculus. For this purpose, we adapt the standard Tait-Martin-Löf technique of parallel reduction.

**Definition 172.** We define a relation  $\rho_{v\Omega} \subseteq (\Lambda_t^\Omega \times \Lambda_t^\Omega) \cup (\Lambda_v^\Omega \times \Lambda_v^\Omega)$ , called parallel  $\beta_{v\Omega}$ -reduction or  $\rho_{v\Omega}$ -reduction, by the following rules:

$$\frac{x \rho_{v\Omega} x}{x \rho_{v\Omega} x} ax \quad \frac{\Omega \rho_{v\Omega} \Omega}{(M)\Omega \rho_{v\Omega} \Omega} ax\Omega \quad \frac{}{(\Omega)M \rho_{v\Omega} \Omega} \Omega_l$$

$$\frac{V \rho_{v\Omega} W}{V! \rho_{v\Omega} W!} !$$

$$\frac{M \rho_{v\Omega} M' \quad V \rho_{v\Omega} V'}{(\lambda x M)!V! \rho_{v\Omega} M'[V'/x]} \rho \quad \frac{M \rho_{v\Omega} M' \quad N \rho_{v\Omega} N'}{MN \rho_{v\Omega} M'N'} @$$

$$\frac{M \rho_{v\Omega} M'}{\lambda x M \rho_{v\Omega} \lambda x M'} \lambda$$

**Remark 173.**

1. The  $\rho_{v\Omega}$ -reduction is reflexive (the proof is by straightforward induction on the expression  $E$ ).
2.  $\beta_{v\Omega} \subseteq \rho_{v\Omega}$  (the proof is by straightforward induction on the length of the derivation of  $E \beta_{v\Omega} E'$ , by exploiting remark 173.1).

**Lemma 174.**  $\rho_{v\Omega} \subseteq \beta_{v\Omega}^*$ .

PROOF. By induction on the length of the derivation of  $E \rho_{v\Omega} E'$ , we prove that  $E \beta_{v\Omega}^* E'$ , for every expressions  $E, E'$ . Let us consider the last rule of this derivation.

If it is the  $ax$ - or  $ax\Omega$ -rule, then we conclude that  $E \beta_{v\Omega}^* E'$  by reflexivity of  $\beta_{v\Omega}^*$ .

If it is the  $x$ -rule with  $x \in \{\Omega_l, \Omega_r, \lambda, !\}$ , then we conclude that  $E \beta_{v\Omega} E'$  by applying the  $x$ -rule.

If it is the  $\rho$ -rule then  $E = (\lambda x M)!V!$  and  $E' = M'[V'/x]$  for some terms  $M, M'$  and values  $V, V'$  such that  $M \rho_{v\Omega} M'$  and  $V \rho_{v\Omega} V'$ . By induction hypothesis,  $M \beta_{v\Omega}^* M'$  and  $V \beta_{v\Omega}^* V'$ , hence  $E \beta_{v\Omega}^* (\lambda x M)!V! \beta_{v\Omega}^* (\lambda x M)!V! \beta_{v\Omega} E'$ .

If it is the  $@$ -rule then  $E = MN$  and  $E' = M'N'$  for some terms  $M, M', N, N'$  such that  $M \rho_{v\Omega} M'$  and  $N \rho_{v\Omega} N'$ . By induction hypothesis,  $M \beta_{v\Omega}^* M'$  and  $N \beta_{v\Omega}^* N'$ , thus  $E \beta_{v\Omega}^* M'N \beta_{v\Omega}^* E'$ .  $\square$

**Corollary 175.**  $\rho_{\nu\Omega}^* = \beta_{\nu\Omega}^*$ .

PROOF.

$\subseteq$ : Proof by straightforward induction on the number  $n \in \mathbb{N}$  of steps of the  $\rho_{\nu\Omega}$ -reduction, by exploiting lemma 174.

$\supseteq$ : The proof is by straightforward induction on the number  $n \in \mathbb{N}$  of steps of the  $\beta_{\nu\Omega}$ -reduction, by exploiting remark 173.2.  $\square$

**Lemma 176.** *For all expressions  $E, E'$ , values  $V, V'$  and variable  $x$ , if  $E \rho_{\nu\Omega} E'$  and  $V \rho_{\nu\Omega} V'$  then  $E[V/x] \rho_{\nu\Omega} E'[V'/x]$ .*

PROOF. By induction on the length of the derivation of  $E \rho_{\nu\Omega} E'$ . Let us consider the last rule of this derivation.

If it is the  $ax$ -rule then  $E = y = E'$ . If  $y = x$  then  $E[V/x] = V$  and  $E'[V'/x] = V'$ , thus  $E \rho_{\nu\Omega} E'$  by hypothesis. Otherwise  $y \neq x$  and so  $E[V/x] = y = E'[V'/x]$ , hence  $E \rho_{\nu\Omega} E'$  by remark 173.1.

If it is the  $ax_{\Omega}$ -rule then  $E = \Omega = E'$  and thus  $E[V/x] = \Omega = E'[V'/x]$ , therefore  $E \rho_{\nu\Omega} E'$  by remark 173.1.

If it is the  $\Omega_l$ - (resp.  $\Omega_r$ -) rule then  $E = \Omega M$  (resp.  $E = M\Omega$ ) and  $E' = \Omega$ , hence  $E[V/x] = (\Omega)M[V/x]$  (resp.  $(M[V/x])\Omega$ ) and  $E' = \Omega$ , therefore  $E \rho_{\nu\Omega} E'$  by the  $\Omega_l$ - (resp.  $\Omega_r$ -) rule.

If it is the  $!$ -rule then  $E = W^!$  and  $E' = W'^!$  for some values  $W, W'$  such that  $W \rho_{\nu\Omega} W'$ , thus  $E[V/x] = (W[V/x])^!$  and  $E'[V'/x] = (W'[V'/x])^!$ . By induction hypothesis  $W[V/x] \rho_{\nu\Omega} W'[V'/x]$ , so  $E[V/x] \rho_{\nu\Omega} E'[V'/x]$  by the  $!$ -rule.

If it is the  $\lambda$ -rule then  $E = \lambda y M$  and  $E' = \lambda y M'$  for some terms  $M, M'$  such that  $M \rho_{\nu\Omega} M'$ , furthermore we can suppose without loss of generality that  $y \neq x$ , thus  $E[V/x] = \lambda y M[V/x]$  and  $E'[V'/x] = \lambda y M'[V'/x]$ . By induction hypothesis  $M[V/x] \rho_{\nu\Omega} M'[V'/x]$ , hence  $E[V/x] \rho_{\nu\Omega} E'[V'/x]$  by the  $\lambda$ -rule.

If it is the  $@$ -rule then  $E = MN$  and  $E' = M'N'$  for some terms  $M, M', N, N'$  such that  $M \rho_{\nu\Omega} M'$  and  $N \rho_{\nu\Omega} N'$ , thus  $E[V/x] = (M[V/x])N[V/x]$  and  $E'[V'/x] = (M'[V'/x])N'[V'/x]$ . By induction hypothesis  $M[V/x] \rho_{\nu\Omega} M'[V'/x]$  and  $N[V/x] \rho_{\nu\Omega} N'[V'/x]$ , hence  $E[V/x] \rho_{\nu\Omega} E'[V'/x]$  by the  $@$ -rule.

If it is the  $\rho$ -rule then  $E = (\lambda y M)^! W^!$  and  $E' = M'[W'/y]$  for some terms  $M, M'$  and values  $W, W'$  such that  $M \rho_{\nu\Omega} M'$  and  $W \rho_{\nu\Omega} W'$ , furthermore we can suppose without loss of generality that  $y \notin \text{fv}(V') \cup \{x\}$ , therefore  $E[V/x] = (\lambda y M[V/x])^! W[V/x]^!$  and  $E'[V'/x] = M'[W'/y][V'/x] = M'[V'/x][W'[V'/x]/y]$  by lemma 99.1. Thus  $M[V/x] \rho_{\nu\Omega} M'[V'/x]$  and  $W[V/x] \rho_{\nu\Omega} W'[V'/x]$  by induction hypothesis, hence  $E[V/x] \rho_{\nu\Omega} E'[V'/x]$  by the  $\rho$ -rule.  $\square$

**Theorem 177** (Strong confluence for  $\rho_{\nu\Omega}$ ). *Let  $E, E_1, E_2$  be expressions: if  $E \rho_{\nu\Omega} E_1$  and  $E \rho_{\nu\Omega} E_2$  with  $E_1 \neq E_2$ , then there exists an expression  $F$  such that  $E_1 \rho_{\nu\Omega} F$  and  $E_2 \rho_{\nu\Omega} F$ .*

PROOF.

□

**Corollary 178** (Confluence for  $\beta_{\sqrt{\Omega}}$ ). *Let  $E, E_1, E_2$  be expressions: if  $E \beta_{\sqrt{\Omega}}^* E_1$  and  $E \beta_{\sqrt{\Omega}}^* E_2$  with  $E_1 \neq E_2$ , then there exists an expression  $F$  such that  $E_1 \beta_{\sqrt{\Omega}}^* F$  and  $E_2 \beta_{\sqrt{\Omega}}^* F$ .*



## Chapter 4

# Two symmetrical call-by-value Krivine abstract machines

Abstract machines play an important role in the implementation of programming languages. The reason abstract machines are so useful is because, on the one hand, they are sufficiently “abstract” to relate easily to other kinds of mathematical semantics, such as equational semantics; on the other hand, they are sufficiently “machine-like” to be easily implementable on real machines.

For the ordinary (call-by-name)  $\lambda$ -calculus, the most remarkable example of abstract machine is the Krivine’s machine (KAM) [Kri85, Kri07, DR04]. For the call-by-value  $\lambda$ -calculus, the first abstract machine was the Landin’s SECD [Lan65], another more recent example is the Leroy’s ZINC [Ler90].

We introduce two versions of the KAM for our call-by-value  $\lambda$ -calculus  $\Lambda_{\text{CBV}}$ , that one without environments (closer to  $\hat{\beta}_v$ -reduction) and that one with environments (closer to what happens in implementations of functional programming languages). Both versions have two subversions: the left-hand one and the right-hand one, which are perfectly symmetric. Our approach is more theoretical and “ $\lambda$ -calculus-like” (as in [Kri85, Kri07, DR04, dC09]) than the abstract machines defined in [Lan65, Ler90].

### 4.1 The versions without environments

**Definition 179** (Stack, process). *A stack is a finite sequence of expressions.*

*A process is a pair  $(M, \pi)$ , denoted by  $M * \pi$ , where  $M$  is a term and  $\pi$  is a stack.*

In other words, a process is a non-empty stack whose first component is a term.

Intuitively, a process can be seen as a program in execution.

**Notation.** Let  $\pi = (E_1, \dots, E_n)$  be a stack with  $n \in \mathbb{N}$ : if  $n = 0$  we denote  $\pi$  by  $\emptyset$ ; for every expression  $E$ , we denote by  $E \cdot \pi$  (resp.  $\pi \cdot E$ ) the stack  $(E, E_1, \dots, E_n)$  (resp.  $(E_1, \dots, E_n, E)$ ); moreover, we denote  $E \cdot \emptyset$  and  $\emptyset \cdot E$  by  $E$ .

**Definition 180.** We define two call-by-value Krivine abstract machines without environments  $\mathsf{K}^l$  (the left CBV-KAM) and  $\mathsf{K}^r$  (the right CBV-KAM) by the following reduction rules:

- this reduction rule is common to  $\mathsf{K}^l$  and  $\mathsf{K}^r$

$$\text{swap} \quad V^! * N \cdot \pi \rightarrow N * V \cdot \pi$$

- these reduction rules are specific for  $\mathsf{K}^l$

$$\begin{array}{l} \text{push}_l \quad (M)N * \pi \rightarrow M * N \cdot \pi \\ \text{pop}_l \quad V^! * \lambda x M \cdot \pi \rightarrow M[V/x] * \pi \end{array}$$

- these reduction rules are specific for  $\mathsf{K}^r$

$$\begin{array}{l} \text{push}_r \quad (M)N * \pi \rightarrow N * M \cdot \pi \\ \text{pop}_r \quad (\lambda x M)^! * V \cdot \pi \rightarrow M[V/x] * \pi \end{array}$$

**Remark 181.** The reduction rules for  $\mathsf{K}^l$  (resp.  $\mathsf{K}^r$ ) are “strongly deterministic” (i.e. they form a partial map from the set of processes to the set of processes): for every process  $M * \pi$  there exists at most one process  $M' * \pi'$  such that  $M * \pi \rightarrow M' * \pi'$  according to a reduction rule of  $\mathsf{K}^l$  (resp.  $\mathsf{K}^r$ ).

Note the different role played by values and terms in  $\mathsf{K}^l$  and  $\mathsf{K}^r$ 's stacks respectively: in  $\mathsf{K}^l$  (resp.  $\mathsf{K}^r$ )'s stack, values have to be seen as functions (resp. arguments) and terms have to be seen as arguments (resp. functions).

The fact that only the  $\text{pop}_c$  rule (with  $c \in \{l, r\}$ ) performs a substitution corresponds to the call-by-value constraint for reduction: the argument in a  $\beta_v$ -redex has to be a value. The  $\text{push}_l$  (resp.  $\text{push}_r$ ) and  $\text{swap}$  rules impose the call-by-value strategy reducing the “leftmost-(resp. rightmost-)outermost”  $\beta_v$ -redex.

**Remark 182.** Let  $M, M'$  be terms,  $E$  be an expression and  $\pi, \pi'$  be stacks: by definition, if  $M * \pi \rightarrow_x M' * \pi'$  with  $x \in \{\text{pop}_l, \text{pop}_r, \text{push}_l, \text{push}_r, \text{swap}\}$ , then  $M * \pi \cdot E \rightarrow_x M' * \pi' \cdot E$  for every expression  $E$ .

**Definition 183.** With every process  $M * \pi$  is associated a term  $\overline{M * \pi}^l$  and a term  $\overline{M * \pi}^r$  defined by induction on the length of  $\pi$  as follows:

$$\begin{array}{ll} \overline{M * \emptyset}^l := M & \overline{M * \emptyset}^r := M \\ \overline{M * V \cdot \pi}^l := \overline{(V^!)M * \pi}^l & \overline{M * V \cdot \pi}^r := \overline{(M)V^! * \pi}^r \\ \overline{M * N \cdot \pi}^l := \overline{(M)N * \pi}^l & \overline{M * N \cdot \pi}^r := \overline{(N)M * \pi}^r \end{array}$$

Roughly speaking, the stack  $\pi$  in a process  $M * \pi$  can be seen as the “applicative closure” of the term  $M$ , and the function  $\overline{(\ )}^1$  (resp.  $\overline{(\ )}^r$ ) allows to rebuild the term corresponding to a given process of  $\mathsf{K}^1$  (resp.  $\mathsf{K}^r$ ), taking into account the swapped application in the stack i.e. the different role played by terms and values in the  $\mathsf{K}^1$  (resp.  $\mathsf{K}^r$ )’s stack.

**Lemma 184.** *Let  $M, N$  be terms,  $V$  be a value and  $\pi$  be a stack.*

1.  $\overline{M * \pi \cdot V}^1 = (V^1) \overline{M * \pi}^1$ .
2.  $\overline{M * \pi \cdot N}^1 = (\overline{M * \pi}^1) N$ .
3.  $\overline{M * \pi \cdot V}^r = (\overline{M * \pi}^r) V^1$ .
4.  $\overline{M * \pi \cdot N}^r = (N) \overline{M * \pi}^r$ .

PROOF. All the proofs are by induction on the length of  $\pi$ .

1. If  $\pi = \emptyset$ , then  $\overline{M * \pi \cdot V}^1 = \overline{M * V}^1 = \overline{(V^1) M * \emptyset}^1 = (V^1) M = (V^1) \overline{M * \emptyset}^1 = (V^1) \overline{M * \pi}^1$ .  
 If  $\pi = W \cdot \pi'$  where  $W$  is a value, then  $\overline{M * \pi \cdot V}^1 = \overline{(W^1) M * \pi' \cdot V}^1 = (V^1) \overline{(W^1) M * \pi'}^1 = (V^1) \overline{M * \pi}^1$  (the central identity holds by induction hypothesis).  
 If  $\pi = L \cdot \pi'$  where  $L$  is a term, then  $\overline{M * \pi \cdot V}^1 = \overline{(M) L * \pi' \cdot V}^1 = (V^1) \overline{(M) L * \pi'}^1 = (V^1) \overline{M * \pi}^1$  (the central identity holds by induction hypothesis).
2. If  $\pi = \emptyset$ , then  $\overline{M * \pi \cdot N}^1 = \overline{M * N}^1 = \overline{(M) N * \emptyset}^1 = (M) N = (\overline{M * \emptyset}^1) N = (\overline{M * \pi}^1) N$ .  
 If  $\pi = W \cdot \pi'$  where  $W$  is a value, then  $\overline{M * \pi \cdot N}^1 = \overline{(W^1) M * \pi' \cdot N}^1 = ((W^1) M * \pi') N = (\overline{M * \pi}^1) N$  (the central identity holds by induction hypothesis).  
 If  $\pi = L \cdot \pi'$  where  $L$  is a term, then  $\overline{M * \pi \cdot N}^1 = \overline{(M) L * \pi' \cdot N}^1 = ((M) L * \pi') N = (\overline{M * \pi}^1) N$  (the central identity holds by induction hypothesis).
3. If  $\pi = \emptyset$ , then  $\overline{M * \pi \cdot V}^r = \overline{M * V}^r = \overline{(M) V^1 * \emptyset}^r = (M) V^1 = (\overline{M * \emptyset}^r) V^1 = (\overline{M * \pi}^r) V^1$ .  
 If  $\pi = W \cdot \pi'$  where  $W$  is a value, then  $\overline{M * \pi \cdot V}^r = \overline{(M) W^1 * \pi' \cdot V}^r = ((M) W^1 * \pi') V^1 = (\overline{M * \pi}^r) V^1$  (the central identity holds by induction hypothesis).  
 If  $\pi = L \cdot \pi'$  where  $L$  is a term, then  $\overline{M * \pi \cdot V}^r = \overline{(L) M * \pi' \cdot V}^r = ((L) M * \pi') V^1 = (\overline{M * \pi}^r) V^1$  (the central identity holds by induction hypothesis).
4. If  $\pi = \emptyset$ , then  $\overline{M * \pi \cdot N}^r = \overline{M * N}^r = \overline{(N) M * \emptyset}^r = (N) M = (N) \overline{M * \emptyset}^r = (N) \overline{M * \pi}^r$ .

If  $\pi = W \cdot \pi'$  where  $W$  is a value, then  $\overline{M * \pi \cdot N}^r = \overline{(M)W^! * \pi' \cdot N}^r = (N)\overline{(M)W^! * \pi'}^r = (N)\overline{M * \pi'}^r$  (the central identity holds by induction hypothesis).

If  $\pi = L \cdot \pi'$  where  $L$  is a term, then  $\overline{M * \pi \cdot N}^r = \overline{(L)M * \pi' \cdot N}^r = (N)\overline{(L)M * \pi'}^r = (N)\overline{M * \pi'}^r$  (the central identity holds by induction hypothesis).  $\square$

Now we compare the two CBV-KAMs with  $\hat{\beta}_v$ -reduction, more precisely we compare  $\mathbf{K}^l$ (resp.  $\mathbf{K}^r$ )'s reduction rules with  $\hat{\beta}_{vl}$ -(resp.  $\hat{\beta}_{vr}$ -)reduction.

**Lemma 185.** *Let  $M$  be a term and  $\pi$  be a stack.*

1. If  $M \hat{\beta}_v M'$  (resp.  $M \hat{\beta}_{vl} M'$ ) then  $\overline{M * \pi}^l \hat{\beta}_v \overline{M' * \pi}^l$  (resp.  $\overline{M * \pi}^l \hat{\beta}_{vl} \overline{M' * \pi}^l$ ).
2. If  $M \hat{\beta}_v M'$  (resp.  $M \hat{\beta}_{vr} M'$ ) then  $\overline{M * \pi}^r \hat{\beta}_v \overline{M' * \pi}^r$  (resp.  $\overline{M * \pi}^r \hat{\beta}_{vr} \overline{M' * \pi}^r$ ).

PROOF. Both proofs are by induction on the length of  $\pi$ .

1. If  $\pi = \emptyset$ , then  $\overline{M * \pi}^l = \overline{M * \emptyset}^l = M \hat{\beta}_v M' = \overline{M' * \emptyset}^l = \overline{M' * \pi}^l$ .  
If  $\pi = V\pi'$  where  $V$  is a value, then  $\overline{M * \pi}^l = \overline{(V^!)M * \pi'}^l \hat{\beta}_v \overline{(V^!)M' * \pi'}^l = \overline{M' * \pi}^l$  (the central relation holds by induction hypothesis, since  $(V^!)M \hat{\beta}_v (V^!)M'$ ).  
If  $\pi = N\pi'$  where  $N$  is a term, then  $\overline{M * \pi}^l = \overline{(M)N * \pi'}^l \hat{\beta}_v \overline{(M')N * \pi'}^l = \overline{M' * \pi}^l$  (the central relation holds by induction hypothesis, since  $(M)N \hat{\beta}_v (M')N$ ).

The proof that  $M \hat{\beta}_{vl} M'$  implies  $\overline{M * \pi}^l \hat{\beta}_{vl} \overline{M' * \pi}^l$  is analogous, it suffices to replace  $\hat{\beta}_v$  by  $\hat{\beta}_{vl}$ .

2. If  $\pi = \emptyset$ , then  $\overline{M * \pi}^r = \overline{M * \emptyset}^r = M \hat{\beta}_v M' = \overline{M' * \emptyset}^r = \overline{M' * \pi}^r$ .  
If  $\pi = V\pi'$  where  $V$  is a value, then  $\overline{M * \pi}^r = \overline{(M)V^! * \pi'}^r \hat{\beta}_v \overline{(M')V^! * \pi'}^r = \overline{M' * \pi}^r$  (the central relation holds by induction hypothesis, since  $(M)V^! \hat{\beta}_v (M')V^!$ ).  
If  $\pi = N\pi'$  where  $N$  is a term, then  $\overline{M * \pi}^r = \overline{(N)M * \pi'}^r \hat{\beta}_v \overline{(N)M' * \pi'}^r = \overline{M' * \pi}^r$  (the central relation holds by induction hypothesis, since  $(N)M \hat{\beta}_v (N)M'$ ).

The proof that  $M \hat{\beta}_{vr} M'$  implies  $\overline{M * \pi}^r \hat{\beta}_{vr} \overline{M' * \pi}^r$  is analogous, it suffices to replace  $\hat{\beta}_v$  by  $\hat{\beta}_{vr}$ .  $\square$

Notice that  $M \hat{\beta}_v M'$  does not entail either  $\overline{M * \pi}^l \hat{\beta}_{vl} \overline{M' * \pi}^l$  or  $\overline{M * \pi}^r \hat{\beta}_{vr} \overline{M' * \pi}^r$ ; for example, take  $M = ((z^!)z^!)((\lambda x x^!)^!)y^!$  (resp.  $M = ((\lambda x x^!)^!)y^!(z^!)z^!$ ),  $M' = ((z^!)z^!)y^!$  (resp.  $M' = (y^!)(z^!)z^!$ ) and  $\pi = \emptyset$ : then  $M \hat{\beta}_v M'$  but  $\overline{M * \emptyset}^l = \overline{M * \emptyset}^r = M$  which is a  $\hat{\beta}_{vl}$ -(resp.  $\hat{\beta}_{vr}$ -)normal form.

**Proposition 186.** *Let  $M$  and  $M'$  be terms, let  $\pi$  and  $\pi'$  be stacks:*

- if  $M * \pi \rightarrow_{\text{pop}_l} M' * \pi'$  (resp.  $M * \pi \rightarrow_{\text{pop}_r} M' * \pi'$ ) then  $\overline{M * \pi}^l \hat{\beta}_{vl} \overline{M' * \pi'}^l$  (resp.  $\overline{M * \pi}^r \hat{\beta}_{vr} \overline{M' * \pi'}^r$ );

- if  $M * \pi \rightarrow_x M' * \pi'$  where  $x \in \{\text{push}_l, \text{swap}\}$  (resp.  $x \in \{\text{push}_r, \text{swap}\}$ ) then  $\overline{M * \pi}^l = \overline{M' * \pi'}^l$  (resp.  $\overline{M * \pi}^r = \overline{M' * \pi'}^r$ ).

PROOF. If  $M * \pi \rightarrow_{\text{pop}_l} M' * \pi'$  (resp.  $M * \pi \rightarrow_{\text{pop}_r} M' * \pi'$ ) then  $M = V^!$ ,  $\pi = \lambda x N \cdot \pi'$  (resp.  $M = (\lambda x N)^!$ ,  $\pi = V \cdot \pi'$ ) and  $M' = N[V/x]$  for some term  $N$  and value  $V$ ; as  $(\lambda x N)^! V^! \hat{\beta}_{vl}$  (resp.  $\hat{\beta}_{vr}$ )  $N[V/x]$ , by lemma 185.1 (resp. 185.2)  $\overline{M * \pi}^l = \overline{(\lambda x N)^! V^! * \pi'}^l \hat{\beta}_{vl} \overline{M' * \pi'}^l$  (resp.  $\overline{M * \pi}^r = \overline{(\lambda x N)^! V^! * \pi'}^r \hat{\beta}_{vr} \overline{M' * \pi'}^r$ ).

If  $M * \pi \rightarrow_{\text{push}_l} M' * \pi'$  (resp.  $M * \pi \rightarrow_{\text{push}_r} M' * \pi'$ ) then  $M = (L)N$  (resp.  $M = (N)L$ ),  $M' = L$  and  $\pi' = N \cdot \pi$  for some terms  $N$  and  $L$ , so  $\overline{M * \pi}^c = \overline{L * N \cdot \pi}^c = \overline{M' * \pi'}^c$  with  $c = l$  (resp.  $c = r$ ).

If  $M * \pi \rightarrow_{\text{swap}} M' * \pi'$  then  $M = V^!$ ,  $M' = N$ ,  $\pi = N \cdot \pi_0$  and  $\pi' = V \cdot \pi_0$  for some term  $N$  and value  $V$ , so  $\overline{M * \pi}^l = \overline{(V^!)N * \pi_0}^l = \overline{M' * \pi'}^l$  (resp.  $\overline{M * \pi}^r = \overline{(N)V^! * \pi_0}^r = \overline{M' * \pi'}^r$ ).  $\square$

Proposition 186 states the soundness of the CBV-KAM  $K^l$  (resp.  $K^r$ )'s reduction rules with respect to  $\hat{\beta}_{vl}$ - (resp.  $\hat{\beta}_{vr}$ -) reduction (and  $\hat{\beta}_v$ -reduction, by remark 117). Indeed the following is an immediate corollary of proposition 186:

**Corollary 187.** *Let  $M, M'$  be terms and  $\pi, \pi'$  be stacks. If  $M * \pi \rightarrow_l M' * \pi'$  (resp.  $M * \pi \rightarrow_r M' * \pi'$ ) then  $\overline{M * \pi}^l \hat{\beta}_{vl}^- \overline{M' * \pi'}^l$  and  $\overline{M * \pi}^r \hat{\beta}_{vr}^- \overline{M' * \pi'}^r$  (resp.  $\overline{M * \pi}^l \hat{\beta}_v^- \overline{M' * \pi'}^l$  and  $\overline{M * \pi}^r \hat{\beta}_v^- \overline{M' * \pi'}^r$ ). In particular, if  $M * \emptyset \rightarrow_l M' * \pi'$  (resp.  $M * \emptyset \rightarrow_r M' * \pi'$ ) then  $M \hat{\beta}_{vl}^- \overline{M' * \pi'}^l$  and  $M \hat{\beta}_v^- \overline{M' * \pi'}^l$  (resp.  $M \hat{\beta}_{vr}^- \overline{M' * \pi'}^r$  and  $M \hat{\beta}_v^- \overline{M' * \pi'}^r$ ).*

**Remark 188.** Let  $M, M', N$  be terms:  $M \hat{\beta}_{vl}$  (resp.  $\hat{\beta}_{vr}$ )  $M'$  does not entail  $N * M \rightarrow_l^* \text{ (resp. } \rightarrow_r^*) N * M'$ . For example, take  $M = (\lambda y y^!)^! z^!$ ,  $M' = z^!$  and  $N = x^!$ :  $M \hat{\beta}_{vl} M'$  and  $M \hat{\beta}_{vr} M'$  but

$$N * M = x^! * (\lambda y y^!)^! z^! \rightarrow_{\text{swap}} (\lambda y y^!)^! z^! * x \left\{ \begin{array}{l} \rightarrow_{\text{push}_l} (\lambda y y^!)^! * z^! \cdot x \rightarrow_{\text{swap}} z^! * \lambda y y^! \cdot x \rightarrow_{\text{pop}_l} \\ \rightarrow_{\text{push}_r} z^! * (\lambda y y^!)^! \cdot x \rightarrow_{\text{swap}} (\lambda y y^!)^! * z \cdot x \rightarrow_{\text{pop}_r} \end{array} \right\} z^! * x \not\rightarrow$$

and every process in the  $K^l$ - (resp.  $K^r$ -) reduction started with  $N * M$  is different from  $N * M'$ .

As a consequence,  $M \hat{\beta}_v M'$  does not imply that there exist a term  $N$  and a stack  $\pi$  such that  $M * \emptyset \rightarrow_l^* \text{ (resp. } \rightarrow_r^*) N * \pi$  and  $M' = \overline{N * \pi}^l$  (resp.  $M' = \overline{N * \pi}^r$ ). For instance, take  $M = ((\lambda x_1 x_1^!)^! y^!) ((\lambda x_2 x_2^!)^! z^!)$  and  $M' = ((\lambda x_1 x_1^!)^! y^!) z^!$  (resp.  $M' = (y^!) ((\lambda x_2 x_2^!)^! z^!)$ ):  $M \hat{\beta}_v M'$  but

$$\begin{aligned} M * \emptyset &\rightarrow_{\text{push}_l} (\lambda x_1 x_1^!)^! y^! * (\lambda x_2 x_2^!)^! z^! \rightarrow_{\text{push}_l} (\lambda x_1 x_1^!)^! * y^! \cdot (\lambda x_2 x_2^!)^! z^! \rightarrow_{\text{swap}} y^! * \lambda x_1 x_1^! \cdot (\lambda x_2 x_2^!)^! z^! \\ &\rightarrow_{\text{pop}_l} y^! * (\lambda x_2 x_2^!)^! z^! \rightarrow_{\text{swap}} (\lambda x_2 x_2^!)^! z^! * y \rightarrow_{\text{push}_l} (\lambda x_2 x_2^!)^! * z^! \cdot y \rightarrow_{\text{swap}} z^! * \lambda x_2 x_2^! \cdot y \rightarrow_{\text{pop}_l} z^! * y \not\rightarrow \\ \text{(resp. } M * \emptyset &\rightarrow_{\text{push}_r} (\lambda x_2 x_2^!)^! z^! * (\lambda x_1 x_1^!)^! y^! \rightarrow_{\text{push}_r} z^! * (\lambda x_2 x_2^!)^! \cdot (\lambda x_1 x_1^!)^! y^! \rightarrow_{\text{swap}} (\lambda x_2 x_2^!)^! * z \cdot (\lambda x_1 x_1^!)^! y^! \\ &\rightarrow_{\text{pop}_r} z^! * (\lambda x_1 x_1^!)^! y^! \rightarrow_{\text{swap}} (\lambda x_1 x_1^!)^! y^! * z \rightarrow_{\text{push}_r} y^! * (\lambda x_1 x_1^!)^! \cdot z \rightarrow_{\text{swap}} (\lambda x_1 x_1^!)^! * y \cdot z \rightarrow_{\text{pop}_r} y^! * z \not\rightarrow \end{aligned}$$

and no process  $N * \pi$  in the  $K^l$ - (resp.  $K^r$ -) reduction started with  $M * \emptyset$  is such that  $M' = \overline{N * \pi}^l$  (resp.  $M' = \overline{N * \pi}^r$ ).

What makes both implications of remark 188 fail is that CBV-KAM  $K^l$  (resp.  $K^r$ )'s reduction rules correspond to the call-by-value strategy reducing the “leftmost- (resp. rightmost-) outermost”  $\beta_v$ -redex, i.e. the  $\hat{\beta}_{vl}$ - (resp.  $\hat{\beta}_{vr}$ -) reduction.

**Proposition 189.** *Let  $M$  and  $M'$  be terms. If  $M \hat{\beta}_{vl}$  (resp.  $\hat{\beta}_{vr}$ )  $M'$  then there exist a term  $N$  and a stack  $\pi$  such that  $M * \emptyset \rightarrow_1^+$  (resp.  $\rightarrow_r^+$ )  $N * \pi$  and  $M' = \overline{N * \pi}^l$  (resp.  $M' = \overline{N * \pi}^r$ ).*

PROOF. By induction on  $M \in \Lambda_t$ . Let us consider the last rule of the derivation of  $M \hat{\beta}_{vl}$  (resp.  $\hat{\beta}_{vr}$ )  $M'$ .

If it is the  $\beta$ -rule, then  $M = (\lambda x N)^! V^!$  and  $M' = N[V/x]$  for some term  $N$  and value  $V$ , hence  $M * \emptyset \rightarrow_{\text{push}_l} (\lambda x N)^! * V^! \rightarrow_{\text{swap}} V^! * \lambda x N \rightarrow_{\text{pop}_l} N[V/x] * \emptyset$  (resp.  $M * \emptyset \rightarrow_{\text{push}_r} V^! * (\lambda x N)^! \rightarrow_{\text{swap}} (\lambda x N)^! * V \rightarrow_{\text{pop}_r} N[V/x] * \emptyset$ ), where  $M' = \overline{N[V/x] * \emptyset}^l$  (resp.  $M' = \overline{N[V/x] * \emptyset}^r$ ).

If it is the  $@_{rv}$ - (resp.  $@_{lv}$ -) rule, then  $M = (V^!)L$  and  $M' = (V^!)L'$  (resp.  $M = (L)V^!$  and  $M' = (L')V^!$ ) for some terms  $L, L'$  and value  $V$  with  $L \hat{\beta}_{vl}$  (resp.  $\hat{\beta}_{vr}$ )  $L'$ , thus there exist a term  $N$  and a stack  $\pi$  such that  $L * \emptyset \rightarrow_1^+$  (resp.  $\rightarrow_r^+$ )  $N * \pi$  and  $L' = \overline{N * \pi}^l$  (resp.  $L' = \overline{N * \pi}^r$ ) by induction hypothesis; hence  $M * \emptyset \rightarrow_{\text{push}_l}$  (resp.  $\rightarrow_{\text{push}_r}$ )  $V^! * L \rightarrow_{\text{swap}} L * V \rightarrow_1^+$  (resp.  $\rightarrow_r^+$ )  $N * \pi \cdot V$  by remark 182, where  $M' = (V^!) \overline{N * \pi}^l = \overline{N * \pi \cdot V}^l$  (resp.  $M' = (\overline{N * \pi}^r) V^! = \overline{N * \pi \cdot V}^r$ ) by lemma 184.1 (resp. 184.3).

If it is the  $@_l$ - (resp.  $@_r$ -) rule, then  $M = (L_2)L_1$  and  $M' = (L'_2)L_1$  (resp.  $M = (L_1)L_2$  and  $M' = (L_1)L'_2$ ) for some terms  $L_1, L_2, L'_2$  with  $L_2 \hat{\beta}_{vl}$  (resp.  $\hat{\beta}_{vr}$ )  $L'_2$ , thus there exist an expression  $N$  and a stack  $\pi$  such that  $L_2 * \emptyset \rightarrow_1^+$  (resp.  $\rightarrow_r^+$ )  $N * \pi$  and  $L'_2 = \overline{N * \pi}^l$  (resp.  $L'_2 = \overline{N * \pi}^r$ ) by induction hypothesis; hence  $M * \emptyset \rightarrow_{\text{push}_l}$  (resp.  $\rightarrow_{\text{push}_r}$ )  $L_2 * L_1 \rightarrow_1^+$  (resp.  $\rightarrow_r^+$ )  $N * \pi \cdot L_1$  by remark 182, where  $M' = (\overline{N * \pi}^l) L_1 = \overline{N * \pi \cdot L_1}^l$  (resp.  $M' = (L_1) \overline{N * \pi}^r = \overline{N * \pi \cdot L_1}^r$ ) by lemma 184.2 (resp. 184.4).  $\square$

Intuitively, proposition 189 is a sort of converse to proposition 186, i.e. it states the “completeness” of the CBV-KAM  $K^l$  (resp.  $K^r$ )'s reduction rules with respect to  $\hat{\beta}_{vl}$ - (resp.  $\hat{\beta}_{vr}$ -) reduction.

**Proposition 190.** *Let  $M$  and  $M'$  be terms. If  $M \hat{\beta}_v M'$  then there exist an expression  $E$  and a stack  $\pi$  such that  $M * \emptyset \rightarrow^+ E * \pi$  and  $M' \hat{\beta}_v^* \overline{E * \pi}$ .*

PROOF. By induction on  $M \in \Lambda_t$ .

If  $M = (\lambda x N)^! V^!$  and  $M' = N[V/x]$ , then  $M * \emptyset \rightarrow_{\text{push}} V^! * (\lambda x N)^! \rightarrow_{\text{access}} V * (\lambda x N)^! \rightarrow_{\text{swap}} (\lambda x N)^! * V \rightarrow_{\text{access}} \lambda x N * V \rightarrow_{\text{pop}} N[V/x] * \emptyset$ , where  $\overline{N[V/x] * \emptyset} = M'$ .

If  $M = (L)N$  and  $M' = (L)N'$  with  $N \hat{\beta}_v N'$ , then by induction hypothesis there exist an expression  $E$  and a stack  $\pi$  such that  $N * \emptyset \rightarrow^+ E * \pi$  and  $N' \hat{\beta}_v^* \overline{E * \pi}$ ; hence  $M * \emptyset \rightarrow_{\text{push}} N * L \rightarrow^+ E * \pi \cdot L$  by remark 182, where  $M' \hat{\beta}_v^* (L)\overline{E * \pi} = \overline{E * \pi} \cdot L$  by lemma 184.1.  $\square$  WHY?

## 4.2 The versions with environments

We recall that the set of variables (resp. values; terms) of  $\Lambda_{\text{CBV}}$  is denoted by  $\mathcal{V}$  (resp.  $\Lambda_v$ ;  $\Lambda_t$ ).

**Definition 191** (Environment). *For every  $p \in \mathbb{N}$ , we define a set  $\mathcal{E}_p$ , by induction on  $p$ , as follows:*

- $\mathcal{E}_0 = \mathcal{V} \rightarrow_{\text{fin}} \emptyset$  (i.e. the set containing only the empty function  $\perp$ );
- if  $p > 0$  then  $\mathcal{E}_p = \mathcal{V} \rightarrow_{\text{fin}} (\Lambda_v \times \mathcal{E}_{p-1})$ .

We set  $\mathcal{E} = \bigcup_{p \in \mathbb{N}} \mathcal{E}_p$ , whose elements are called environments. For every  $e \in \mathcal{E}$ , we denote by  $d(e)$  the least  $p \in \mathbb{N}$  such that  $e \in \mathcal{E}_p$ .

Intuitively, an environment can be seen as a kind of heap memory used for dynamic memory allocation.

**Remark 192.** For every  $p \in \mathbb{N}$ ,  $\mathcal{E}_p \subseteq \mathcal{E}_{p+1}$ . The proof is a straightforward induction on  $p \in \mathbb{N}$ . The empty function is in  $\mathcal{E}_1$ , so  $\mathcal{E}_0 \subseteq \mathcal{E}_1$ . Let  $p > 0$  and  $e \in \mathcal{E}_p$ : if  $\text{dom}(e) = \{x_1, \dots, x_n\}$  for some  $n \in \mathbb{N}$ , then for every  $1 \leq i \leq n$  there exist a value  $V_i$  and  $e_i \in \mathcal{E}_{p-1}$  such that  $e(x_i) = (V_i, e_i)$ ; by induction hypothesis,  $\mathcal{E}_{p-1} \subseteq \mathcal{E}_p$  and thus  $e_i \in \mathcal{E}_p$ , hence  $e(x_i) \in \Lambda_v \times \mathcal{E}_p$ ; therefore  $e \in \mathcal{V} \rightarrow_{\text{fin}} (\Lambda_v \times \mathcal{E}_p) = \mathcal{E}_{p+1}$ , whence  $\mathcal{E}_p \subseteq \mathcal{E}_{p+1}$ .

**Definition 193** (Closure). *The set  $\mathcal{C}_v$  of value (resp. term) closures is defined by  $\mathcal{C}_v = \Lambda_v \times \mathcal{E}$  (resp.  $\mathcal{C}_t = \Lambda_t \times \mathcal{E}$ ). The set  $\mathcal{C}$  of closures is defined by  $\mathcal{C} = \mathcal{C}_v \cup \mathcal{C}_t$ .*

*Given  $v = (V, e) \in \mathcal{C}_v$ , we define  $\bar{v} = V[e] \in \Lambda_v$  by induction on  $d(e) \in \mathbb{N}$ :*

- if  $d(e) = 0$  then  $V[e] = V$ ;
- if  $d(e) > 0$  then  $V[e] = V[\overline{e(x_1)}/x_1, \dots, \overline{e(x_n)}/x_n]$  where  $\text{dom}(e) = \{x_1, \dots, x_n\}$  for some  $n \in \mathbb{N}$ .

*Given  $t = (M, e) \in \mathcal{C}_t$ , we define  $\bar{t} = M[e] \in \Lambda_t$  by induction on  $M \in \Lambda_t$ :*

- if  $M = V^!$  for some value  $V$ , then  $M[e] = (V[e])^!$ ;
- if  $M = NL$  for some terms  $N$  and  $L$ , then  $M[e] = (N[e])L[e]$ .

**Remark 194.**

1. With reference to notations used in definition 193, note that  $\bar{v}$  is well-defined for  $v = (V, e) \in \mathcal{C}_v$ : indeed, in the case  $\mathbf{d}(e) > 0$ , for every  $1 \leq i \leq n$  there exists a value  $V_i$  and  $e_i \in \mathcal{E}_{\mathbf{d}(e)-1}$  such that  $e(x_i) = (V_i, e_i)$ , hence  $\mathbf{d}(e_i) \leq \mathbf{d}(e) - 1$  and so  $\overline{e(x_i)}$  is defined by induction hypothesis. Furthermore if  $c \in \mathcal{C}$  is a value (resp. term) closure, then  $\bar{c}$  is a value (resp. term).
2. By definition,  $\mathcal{E} = \mathcal{V} \multimap_{\text{fin}} \mathcal{C}_v$ , i.e. environments are the partial functions with finite domains from the set of variables to the set of value closure.
3.  $\mathcal{C} = \mathcal{C}_v \uplus \mathcal{C}_t$ , since  $\Lambda_v \cap \Lambda_t = \emptyset$ .

Let  $E$  be an expression and  $e$  be an environment: each pair  $(x, v)$  (where  $x$  is a variable and  $v = (V, e')$  is a value closure) in the graph of  $e$  can be seen as a sort of “recursive” explicit substitution in the expression  $E[e]$ , associating  $V[e']$  with the free occurrences of  $x$  in  $E$ .

**Definition 195** (Stack, state). *A stack is a finite sequence of closures.*

A state is a pair  $(t, \pi)$ , denoted by  $t * \pi$ , where  $t$  is a term closure and  $\pi$  is a stack. If  $s = t * (c_1, \dots, c_n)$  for some  $n \in \mathbb{N}$  is a state, then  $\bar{s}$  denotes the term  $(\bar{t})\bar{c}_1 \cdots \bar{c}_n$ .

In other words, a state is a non empty stack whose first component is a term closure.

Intuitively, a state is a program in execution, taking into account also the environment of this execution.

**Definition 196** (Variable convention). *For every value closure  $v = (V, e)$ , we define, by induction on  $\mathbf{d}(e)$ , what means that the value closure  $v$  respects the variable convention;  $v$  respects the variable convention if the following conditions are fulfilled:*

- every bound variable in  $E$  is bound in  $E$  at most once;
- for every bound variable  $x$  in  $E$ ,  $x \notin \text{dom}(e)$ ;
- for every  $v \in \text{im}(e)$ ,  $v$  respects the variable convention.

For every term closure  $t = (M, e)$ , we define, by induction on  $M$ , what means that the term closure  $t$  respects the variable convention:

- if  $M = V^!$  for some value  $V$  then  $t$  respects the variable convention if  $(V, e)$  respects the variable convention;
- if  $M = NL$  for some terms  $N$  and  $L$  then  $t$  respects the variable convention if  $(N, e)$  and  $(L, e)$  respect the variable convention.

We say that state  $t * (c_1, \dots, c_n)$  (where  $n \in \mathbb{N}$ ,  $t$  is a term closure and  $c_i$  is a closure for any  $1 \leq i \leq n$ ) respects the variable convention if the closures  $t, c_1, \dots, c_n$  respect the variable convention.



For instance,  $(\lambda y(\lambda x(x^!)x^!)^!, \perp)$  respects the variable convention, whereas  $(\lambda x(\lambda x x^!)^!, \perp)$  does not.

**Definition 197.** We define two call-by-value Krivine abstract machines with environments  $K_{\text{env}}^l$  (the left CBV-KAM<sub>env</sub>) and  $K_{\text{env}}^r$  (the right CBV-KAM<sub>env</sub>) by the following reduction rules:

- these reductions rule are common to  $K_{\text{env}}^l$  and  $K_{\text{env}}^r$

$$\begin{array}{l} \text{swap} \quad (V^!, e) * (M, e') \cdot \pi \rightarrow (M, e') * (V, e) \cdot \pi \quad \text{if either } V \notin \mathcal{V}, \text{ or } V \in \mathcal{V} \text{ and } V \notin \text{dom}(e) \\ \text{sub} \quad \quad (x^!, e) * \pi \rightarrow (V^!, e') * \pi \quad \quad \text{if } x \in \text{dom}(e) \text{ and } e(x) = (V, e') \end{array}$$

- these reduction rules are specific for  $K_{\text{env}}^l$

$$\begin{array}{l} \text{push}_l \quad \quad (MN, e) * \pi \rightarrow (M, e) * (N, e) \cdot \pi \\ \text{pop}_l \quad (V^!, e') * (\lambda x M, e) \cdot \pi \rightarrow (M, e \cup \{x \mapsto (V, e')\}) * \pi \end{array}$$

- these reduction rules are specific for  $K_{\text{env}}^r$

$$\begin{array}{l} \text{push}_r \quad \quad (MN, e) * \pi \rightarrow (N, e) * (M, e) \cdot \pi \\ \text{pop}_r \quad ((\lambda x M)^!, e) * (V, e') \cdot \pi \rightarrow (M, e \cup \{x \mapsto (V, e')\}) * \pi \end{array}$$

**Remark 198.** The reduction rules for  $K_{\text{env}}^l$  (resp.  $K_{\text{env}}^r$ ) are “strongly deterministic” (i.e. they form a partial map from the set of states to the set of states): for every state  $t * \pi$  there exists at most one state  $t' * \pi'$  such that  $t * \pi \rightarrow t' * \pi'$  according to a reduction rule of  $K_{\text{env}}^l$  (resp.  $K_{\text{env}}^r$ ).



## Chapter 5

# Translations

### 5.1 The typed $\Lambda_{\text{CBV}}$ and boring translations in Linear Logic

A type system is a class of formulas in some language, the purpose of which is to express some properties of  $\lambda$ -terms. By introducing such formulas, as comments in the terms, we construct what we call typed terms, which correspond to programs in a high level programming language. The main connective in these formulas is “ $\rightarrow$ ”, the type  $A \rightarrow B$  being that of the “functions” from  $A$  to  $B$ , that is to say from the set of terms of type  $A$  to the set of terms of type  $B$ .

By a variable declaration, we mean an ordered pair  $(x, A)$ , where  $x$  is a variable of the  $\lambda$ -calculus, and  $A$  is a type. It will be denoted by  $x : A$  instead of  $(x, A)$ . A context  $\Gamma$  is a mapping from a finite set of variables to the set of all types. Thus it is a finite set  $\{x_1 : A_1, \dots, x_k : A_k\}$  of variable declarations, where  $x_1, \dots, x_k$  are distinct variables ; we will denote it by  $x_1 : A_1, \dots, x_k : A_k$  (without the braces). So, in such an expression, the order does not matter. We will say that  $x_i$  is declared of type  $A_i$  in the context  $\Gamma$ . The integer  $k$  may be 0; in that case, we have the empty context.

We will write  $\Gamma, x : A$  in order to denote the context obtained by adding the declaration  $x : A$  to the context  $\Gamma$ , provided that  $x$  is not already declared in  $\Gamma$ .

Given a  $\lambda$ -term  $t$ , a type  $A$ , and a context  $\Gamma$ , we define, by means of the following rules, the notion:  $t$  is of type  $A$  in the context  $\Gamma$  (we will also say : “ $t$  may be given type  $A$  in the context  $\Gamma$ ”); this will be denoted by  $\Gamma \vdash_{\mathcal{L}} t : A$  (or  $\Gamma \vdash t : A$  if there is no ambiguity) :

$$\frac{}{\Gamma, x : A \vdash x : A} \text{ax} \qquad \frac{\Gamma \vdash V : A}{\Gamma \vdash V! : A} !$$
$$\frac{\Gamma, x : A \vdash M : B}{\Gamma \vdash \lambda x M : A \rightarrow B} \rightarrow_i \qquad \frac{\Gamma \vdash M : A \rightarrow B \quad \Gamma \vdash N : A}{\Gamma \vdash MN : B} \rightarrow_e$$

There are two ways to traduce the intuitionistic arrow  $A \rightarrow B$  in Linear Logic with a “call-by-value” style (see [Gir87]).

$$\begin{array}{ll}
X^\circ & := !X & X^\bullet & := X \\
(A \rightarrow B)^\circ & := !(A^\circ \multimap B^\circ) & (A \rightarrow B)^\bullet & := (!A^\bullet \multimap !B^\bullet) \\
(\Gamma \vdash V : A)^\circ & := !\Gamma^\circ \vdash A^\circ & (\Gamma \vdash V : A)^\bullet & := !\Gamma^\bullet \vdash V : A^\bullet \\
(\Gamma \vdash M : A)^\circ & := !\Gamma^\circ \vdash A^\circ & (\Gamma \vdash M : A)^\bullet & := !\Gamma^\bullet \vdash M : !A^\bullet
\end{array}$$

In some sense, the following proposition means that these two translations are equivalent.

**Proposition 199.** *For every formula  $A$  of the implicative fragment,  $A^\circ = !A^\bullet$ .*

PROOF. By induction on the formula  $A$ .

- If  $A = X$ , then  $A^\circ = !X = !A^\bullet$ .
- If  $A = B \rightarrow C$ , then  $B^\circ = !B^\bullet$  and  $C^\circ = !C^\bullet$  by induction hypothesis, so  $A^\circ = !(B^\circ \multimap C^\circ) = !(!B^\bullet \multimap !C^\bullet) = !A^\bullet$ . □

## 5.2 $\sigma_v$ -equivalence

In the ordinary (call-by-name)  $\lambda$ -calculus the  $\sigma$ -equivalence (introduced by Regnier in [Reg92, Reg94]) identifies terms that differ only in their sequential structure (e.g.  $(\lambda x_1 \lambda x_2 u)v_1 v_2$  and  $(\lambda x_2 \lambda x_1 u)v_2 v_1$ ):  $\lambda$ -terms contain pieces of information, which are unnecessary from the operational view-point. The same phenomenon may be found in the call-by-value  $\lambda$ -calculus. So two questions naturally arise for  $\Lambda_{\text{CBV}}$ : find the  $\sigma_v$ -equivalence for  $\Lambda_{\text{CBV}}$ ; find some parallel syntax which identifies  $\sigma_v$ -equivalent terms. In the ordinary  $\lambda$ -calculus, these two questions are answered by means of the Girard’s translation of intuitionistic logic into Linear Logic proof-nets:  $(A \rightarrow B) \rightsquigarrow (!A \multimap B)$ . We give an analogous answer for  $\Lambda_{\text{CBV}}$  by means of the “boring” translation of intuitionistic logic into Linear Logic proof-nets:  $(A \rightarrow B) \rightsquigarrow (!A \multimap !B)$ .

Interestingly, this new  $\sigma_v$ -equivalence relation is not included in the  $\beta_v$ -equivalence, i.e. the  $\sigma_v$ -equivalence identifies distinct  $\beta_v$ -normal terms. We eventually show that two terms are equivalent iff they are translated as the same Linear Logic proof-net.

The  $\sigma_v$ -equivalence is generated by the following rules:

$$\begin{array}{l}
\sigma_1: (\lambda x M)^\dagger N L \simeq (\lambda x M L)^\dagger N \text{ with } x \notin \text{fv}(L); \\
\sigma_2: (\lambda x (\lambda y L)^\dagger)^\dagger M N \simeq (\lambda y (\lambda x L)^\dagger)^\dagger N M; \\
\sigma_3: (M)^\dagger ((\lambda x L)^\dagger)^\dagger N \simeq (\lambda x M L)^\dagger N \text{ with } x \notin \text{fv}(M)
\end{array}$$

None of these rules are included in the  $\beta_v$ -equivalence differently from the standard (call-by-name)  $\lambda$ -calculus, where the  $\sigma$ -equivalence is included in the  $\beta$ -equivalence. In some sense, the  $\beta_v$ -equivalence is incomplete, and the  $\sigma_v$ -equivalence is its completion.

**Theorem 200.** *For every expression  $E$  and  $F$ ,  $E \simeq F$  iff  $E^\bullet = F^\bullet$ .*

## 5.3 CPS

A more significant way than the forgetful functor to embedding the call-by-value  $\lambda$ -calculus  $\Lambda_{\text{CBV}}$  into the ordinary (call-by-name)  $\lambda$ -calculus  $\Lambda$  is the continuation-passing style (CPS) translation. We present two CPS translations, the left one  $()^l$  (already used in [Plo75, Sel01]) and the right one  $()^r$ .

**Definition 201.** *Let  $E$  be an expression.*

*We define, by induction on  $E$ , the right CPS translation of  $E$ , denoted by  $E^r \in \Lambda$ , as follows:*

- $x^r = x$ ;
- $(\lambda x M)^r = \lambda x M^r$ ;
- $(V^l)^r = \lambda k(k)V^r$  with  $k \notin \text{fv}(V)$ ;
- $(MN)^r = \lambda k(N^r)\lambda n(M^r)\lambda m(m)nk$  with  $k, m, n \notin \text{fv}(MN)$ .

*We define, by induction on  $E$ , the left CPS translation of  $E$ , denoted by  $E^l \in \Lambda$ , as follows:*

- $x^l = x$ ;
- $(\lambda x M)^l = \lambda x M^l$ ;
- $(V^l)^l = \lambda k(k)V^l$  with  $k \notin \text{fv}(V)$ ;
- $(MN)^l = \lambda k(M^l)\lambda m(N^l)\lambda n(m)nk$  with  $k, m, n \notin \text{fv}(MN)$ .

Note that the only difference between left and right CPS translations is in the applicative case.

**Remark 202.** For every expression  $E$ ,  $\text{fv}(E^c) = \text{fv}(E)$  with  $c \in \{l, r\}$  (the proof is a straightforward induction on  $E \in \Lambda_{\text{CBV}}$ ).

**Lemma 203** (Substitution). *For every expression  $E$ , value  $V$  and variable  $x$ ,  $(E[V/x])^c = E^c[V^c/x]$  with  $c \in \{l, r\}$ .*

PROOF. By induction on  $E \in \Lambda_{\text{CBV}}$ . Let  $c \in \{l, r\}$ .

If  $E = x$ , then  $E[V/x] = V$  and  $E^c = x$ , so  $(E[V/x])^c = V^c = E^c[V^c/x]$ .

If  $E = y$  for some variable  $y \neq x$ , then  $E[V/x] = y$  and  $E^c = y$ , so  $(E[V/x])^c = y = E^c[V^c/x]$ .

If  $E = \lambda y M$  for some term  $M$ , then we can suppose without loss of generality  $y \notin \text{fv}(V) \cup \{x\}$  (by  $\alpha$ -equivalence), thus  $E[V/x] = \lambda y M[V/x]$  and  $E^c = \lambda y M^c$  with  $(M[V/x])^c = M^c[V^c/x]$  by induction hypothesis, so  $(E[V/x])^c = (\lambda y M[V/x])^c = \lambda y (M[V/x])^c = \lambda y M^c[V^c/x] = E^c[V^c/x]$  since  $y \notin \text{fv}(V^c) \cup \{x\}$  by remark 202.

If  $E = W^!$  for some value  $W$ , then  $E[V/x] = (W[V/x])^!$  and  $E^c = \lambda k(k)W^c$  with  $k \notin \text{fv}(W) \cup \{x\}$  (by  $\alpha$ -equivalence), thus  $k \notin \{x\} \cup \text{fv}(W^c)$  by remark 202, moreover  $(W[V/x])^c = W^c[V^c/x]$  by induction hypothesis, so  $(E[V/x])^c = \lambda k(k)(W[V/x])^c = \lambda k(k)W^c[V^c/x] = E^c[V^c/x]$ .

If  $E = MN$  for some terms  $M$  and  $N$ , then  $E[V/x] = M[V/x]N[V/x]$  and

$$\begin{aligned} E^l &= \lambda k(M^l)\lambda m(N^l)\lambda n(m)nk \\ E^r &= \lambda k(N^r)\lambda n(M^r)\lambda m(m)nk \end{aligned}$$

with  $k, m, n \notin \text{fv}(M) \cup \text{fv}(N) \cup \{x\} = \text{fv}(M^c) \cup \text{fv}(N^c) \cup \{x\}$  by remark 202 and  $\alpha$ -equivalence, moreover  $(M[V/x])^c = M^c[V^c/x]$  and  $(N[V/x])^c = N^c[V^c/x]$  by induction hypothesis, hence

$$\begin{aligned} (E[V/x])^l &= \lambda k(M[V/x])^l \lambda m(N[V/x])^l \lambda n(m)nk = \lambda k(M^l[V^l/x]) \lambda m(N^l[V^l/x]) \lambda n(m)nk = E^l[V^l/x] \\ (E[V/x])^r &= \lambda k(N[V/x])^r \lambda n(M[V/x])^r \lambda m(m)nk = \lambda k(N^r[V^r/x]) \lambda n(M^r[V^r/x]) \lambda m(m)nk = E^r[V^r/x] \end{aligned}$$

□

**Remark 204.** For every values  $V_1$  and  $V_2$ , if  $c \in \{l, r\}$  then  $(V_1^! V_2^!)^c \beta^+ \lambda k(V_1^c) V_2^c k$  with  $k \notin \text{fv}(V_1^c V_2^c)$ . Indeed, let  $k, m, n \notin \text{fv}(V_1^c) \cup \text{fv}(V_2^c) = \text{fv}(V_1^c V_2^c)$  with  $c \in \{l, r\}$ :

$$\begin{aligned} (V_1^! V_2^!)^l &= \lambda k(\lambda k_1(k_1) V_1^!) \lambda m(\lambda k_2(k_2) V_2^!) \lambda n(m)nk \quad \beta \quad \lambda k(\lambda k_1(k_1) V_1^!) \lambda m(\lambda n(m)nk) V_2^! \\ &\quad \beta \quad \lambda k(\lambda m(\lambda n(m)nk) V_2^!) V_1^! \quad \beta \quad \lambda k(\lambda n(V_1^!)nk) V_2^! \quad \beta \quad \lambda k V_1^! V_2^! k. \\ (V_1^! V_2^!)^r &= \lambda k(\lambda k_2(k_2) V_2^r) \lambda n(\lambda k_1(k_1) V_1^r) \lambda m(m)nk \quad \beta \quad \lambda k(\lambda k_2(k_2) V_2^r) \lambda n(\lambda m(m)nk) V_1^r \\ &\quad \beta \quad \lambda k(\lambda n(\lambda m(m)nk) V_1^r) V_2^r \quad \beta \quad \lambda k(\lambda n(V_1^r)nk) V_2^r \quad \beta \quad \lambda k V_1^r V_2^r k. \end{aligned}$$

The following proposition claims that one step of  $\beta_v$ - (and so  $\hat{\beta}_v$ -) reduction is simulated by at least one step of  $\beta_\eta$ -reduction in ordinary  $\lambda$ -calculus, modulo left or right CPS translation.

**Proposition 205.** *Let  $E, E' \in \Lambda_{\text{CBV}}$ : if  $E \beta_v E'$  then  $E^c \beta_\eta^+ E'^c$  with  $c \in \{l, r\}$ .*

PROOF. By induction on  $E \in \Lambda_{\text{CBV}}$ . Let us consider the last rule of the derivation of  $E \beta_{\mathbf{v}} E'$ .

If it is the  $\beta$ -rule, then  $E = (\lambda x M)^! V^!$  and  $E' = M[V/x]$  for some term  $M$  and value  $V$ , hence  $E^c \beta^+ \lambda k(\lambda x M^c) V^c k$  by remark 204 and  $E'^c = (M[V/x])^c = M^c[V^c/x]$  by lemma 203; thus  $E^c \beta^+ \lambda k(\lambda x M^c) V^c k \beta \lambda k(M^c[V^c/x]) k = \lambda k(E'^c) k \eta E'^c$  since  $k \notin \text{fv}(M^c) \cup \text{fv}(V^c) \cup \{x\}$ .

If it is the  $@_l$ -rule, then  $E = MN$  and  $E' = M'N$  where  $M, M', N$  are terms with  $M \beta_{\mathbf{v}} M'$ , hence  $M^c \beta_{\eta^+} M'^c$  with  $c \in \{l, r\}$  by induction hypothesis, so

$$\begin{aligned} E^l &= \lambda k(M^l) \lambda m(N^l) \lambda n(m) n k \beta_{\eta^+} \lambda k(M'^l) \lambda m(N^l) \lambda n(m) n k = E'^l \\ E^r &= \lambda k(N^r) \lambda n(M^r) \lambda m(m) n k \beta_{\eta^+} \lambda k(N^r) \lambda n(M'^r) \lambda m(m) n k = E'^r \end{aligned}$$

since  $\beta_{\eta}$ -reduction passes to context.

If it is the  $@_r$ -rule, then  $E = MN$  and  $E' = MN'$  where  $M, N, N'$  are terms with  $N \beta_{\mathbf{v}} N'$ , hence  $N^c \beta_{\eta^+} N'^c$  with  $c \in \{l, r\}$  by induction hypothesis, so

$$\begin{aligned} E^l &= \lambda k(M^l) \lambda m(N^l) \lambda n(m) n k \beta_{\eta^+} \lambda k(M^l) \lambda m(N'^l) \lambda n(m) n k = E'^l \\ E^r &= \lambda k(N^r) \lambda n(M^r) \lambda m(m) n k \beta_{\eta^+} \lambda k(N'^r) \lambda n(M^r) \lambda m(m) n k = E'^r \end{aligned}$$

since  $\beta_{\eta}$ -reduction passes to context.

If it is the  $\lambda$ -rule, then  $E = \lambda x M$  and  $E' = \lambda x M'$  where  $M$  is a term with  $M \beta_{\mathbf{v}} M'$ , hence for  $c \in \{l, r\}$ ,  $M^c \beta_{\eta^+} M'^c$  by induction hypothesis, so  $E^c = \lambda x M^c \beta_{\eta^+} \lambda x M'^c = E'^c$  since  $\beta_{\eta}$ -reduction passes to context.

If it is the  $!$ -rule, then  $E = V^!$  and  $E' = V'^!$  where  $V, V'$  are values with  $V \beta_{\mathbf{v}} V'$ , hence for  $c \in \{l, r\}$ ,  $V^c \beta_{\eta^+} V'^c$  by induction hypothesis, so  $E^c = \lambda k(k) V^c \beta_{\eta^+} \lambda k(k) V'^c = E'^c$  since  $\beta_{\eta}$ -reduction passes to context.  $\square$





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