# Burden in Henselian Valued Fields 

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#### Abstract

In the spirit of the Ax-Kochen-Ershov principle, we show that in certain cases the burden of a Henselian valued field can be computed in terms of the burden of its residue field and that of its value group. To do so, we first see that the burden of such a field is equal to the burden of its leading term structure. These results are generalisations of Chernikov and Simon's work in [11].


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[^0]
## Introduction

Since the work of Shelah [32], model theorists have defined and studied combinatorial configurations that first order theories may encode, in order to measure their relative tameness. Theories which are able to encode complex configurations are considered less tame. Arising from this classification a complex hierarchy of first order theories. The most important class is undoubtedly that of stable theories; the classes of NIP (Non-Independance Property), simple and NTP ${ }_{2}$ theories (Non-Tree Property of the second kind) are generalisations of stability which have been extensively studied in the last 20 years. The abstract study of these classes leads to a better understanding of algebraic structures, as some algebraic phenomena may or not occur, depending of the complexity of the theory (e.g. [23]). Locating a given concrete first-order theory in this hierarchy is often interesting and challenging, and many methods are known: one can show stability (resp. NIP) by counting types over small sets (resp. coheirs over models). Simple theories are exactly theories with some abstract independence relation ([25]). The situation seems to be different when we look at quantitative versions of these notions, such as the burden. This is a notion of dimension for $\mathrm{NTP}_{2}$ theories, and in order to compute it, a concrete understanding of formulas is required.

Since Ax, Kochen and Ershov's work, Henselian valued fields appear to be perfect playgrounds for this kind of consideration, and many model theoretic questions on valued fields have been reduced to the residue field and value group. The theorem of Delon [12] is one of the first instances: a Henselian valued field of equicharacteristic 0 is NIP if and only if both of its residue field and its value group are NIP ${ }^{1}$. A more quantitative transfer in NIP Henselian valued fields was then showed by Shelah in [33]: a Henselian valued field of equicharacteristic 0 is strongly dependant if and only if both its residue field and its value group are strongly dependant. Both results were generalised to $\mathrm{NTP}_{2}$ theories by Chernikov in [8]: a Henselian valued field of equicharacteristic 0 is $\mathrm{NTP}_{2}$ (resp. strong, of finite burden) if and only if both of its residue field and its value group are $\mathrm{NTP}_{2}$ (resp. strong, of finite burden). This approach is based on the faith that the study of the residue field and the value group might be enough to classify some rather nice valued fields. In the case of equicharacteristic 0 Henselian valued fields, this faith is justified by the theorem of Pas: a Henselian valued field of equicharacteristic zero equipped with an angular component (also called ac-map), eliminates quantifiers relative to the value group and the residue field. It is indeed an important tool for producing transfer principles. However, the theorem of Pas has some limits: considering an ac-map has the disadvantage of adding new definable sets to the structure of valued fields. From the point of view of complexity, it has an impact: any ultraproduct of p-adic fields over a non-principal utrafilter on prime numbers is inp-minimal, i.e. of burden 1 (see [11]), but it is of burden 2 when endowed with an ac-map. Another approach initiated by Basarab and Kuhlmann, is to consider another interpretable sort capturing information from both the value group and the residue field. This lead to the definition of the leading term structure (see Paragraph 1.2.1), also called the RV-sort. Unlike the ac-map, it is always interpretable in the standard languages of valued fields, and adding it to the language does not add definable sets. The study of valued fields together with its RV-sort offers an additional point of view: let us cite Hrushovski and Kazdan's work in motivic integration, where RV-sort structures are used, as opposed to Denef, Loeser and Cluckers' work, where ac-maps are used. More relevant to this paper, Chernikov and Simon prove in [11] that, under some hypothesis, an equicharacteristic zero Henselian valued field is inp-minimal if only if both its residue field and its value group are inp-minimal, going via an intermediate step: they first reduce the question to the RV-sort.

The main aim of this paper is to give a general transfer principle for burden in certain Henselian valued fields. In particular, we provide a full answer to [11, Problems $4.3 \& 4.4]$. Here is an overview of the paper:
The first section consists of preliminaries on pure model theory and on model theory of valued fields

[^1]and groups. We define the burden of a theory, indiscernible sequences and few related lemmas. Then we define the sort RV and recall some relative quantifier elimination results.

Section 2 is dedicated to the proof of the following theorem:
Theorem 2.2. Consider an $\{A\}-\{C\}$-enrichment of a pure exact sequence $\mathcal{M}$ of abelian groups

$$
0 \longrightarrow A \xrightarrow{\iota} B \xrightarrow{\nu} C \longrightarrow 0,
$$

in a language L. Let $\mathcal{D}=\mathcal{D}(x)$ be the set of formulas in the pure language of groups which are conjunction of formulas of the form $\exists y n x=m y$ for $n, m \in \mathbb{N}$. For $D(x) \in \mathcal{D}$ and $A$ an abelian group, $D(A)$ is a subgroup of $A$, and we have

$$
\operatorname{bdn} \mathcal{M}=\max _{D \in \mathcal{D}}(\operatorname{bdn}(A / D(A))+\operatorname{bdn}(D(C)))
$$

In particular:

- If $A / n A$ is finite for all $n \geq 1$, then

$$
\operatorname{bdn} \mathcal{M}=\max _{k \in \mathbb{N}}\left(\operatorname{bdn}(k A)+\operatorname{bdn}\left(C_{[k]}\right)\right)
$$

where $C_{[k]}:=\{c \in C \mid k c=0\}$ is the subgroup of $k$-torsion.

- If $C$ has finite $k$-torsion of all $k \geq 1$, then

$$
\operatorname{bdn} \mathcal{M}=\max _{n \in \mathbb{N}}(\operatorname{bdn}(A / n A)+\operatorname{bdn}(n C))
$$

- If $C$ has finite $n$-torsion and $A / n A$ is finite for all $n \geq 1$, then

$$
\operatorname{bdn} \mathcal{M}=\max (\operatorname{bdn}(A), \operatorname{bdn}(C))
$$

We also give an easy generalisation, notably for short exact sequences of ordered abelian groups.

In Section 3, we prove our main theorems. The first concerns the following partial theories of Henselian valued fields, that we gather here under the name of benign theories of valued fields:

1. Henselian valued fields of characteristic $(0,0)$,
2. algebraically closed valued fields,
3. algebraically maximal Kaplansky valued fields.

Theorem 3.12. Let $\mathcal{K}=(K, \Gamma, k)$ be a benign Henselian valued field, with value group $\Gamma$ and residue field $k$. Then:

$$
\operatorname{bdn}(\mathcal{K})=\max _{n \geq 0}\left(\operatorname{bdn}\left(k^{\star} / k^{\star n}\right)+\operatorname{bdn}(n \Gamma)\right)
$$

The second concerns unramified mixed characteristic Henselian valued fields with perfect residue field.

Theorem 3.21. Let $\mathcal{K}=(K, k, \Gamma)$ be an unramified mixed characteristic Henselian valued field with value group $\Gamma$ and residue field $k$. We denote by $\mathcal{K}_{\mathrm{ac}}^{<\omega}=\left(K,k, \Gamma, \textrm{ac}_{n}, n<\omega\right)$ the structure $\mathcal{K}$ endowed with compatible ac-maps. Assume that the residue field $k$ is perfect. One has

$$
\operatorname{bdn}(\mathcal{K})=\operatorname{bdn}\left(\mathcal{K}_{\mathrm{ac}<\omega}\right)=\max \left(\aleph_{0} \cdot \operatorname{bdn}(k), \operatorname{bdn}(\Gamma)\right)
$$

For both proofs, and as in [11], we proceed in two steps: we show first that the burden of these Henselian valued fields is equal to the burden of their RV-sorts (see Theorem 3.2) and use Section 2 to conclude.

Let us conclude this introduction with some generalities about transfer principles. We summarise the strategy presented above by formalising the reduction in valued fields and pure short exact sequences of abelian groups. We introduce reduction diagrams. It is nothing else than a concise way to picture relative quantifier elimination and by extension, the strategy for proving reduction principles.

Heuristic 0.1. A reduction diagram of a structure $\mathcal{M}$ is a rooted tree such that:

- all nodes are pure sorts of $\mathcal{M}$ (in some $\emptyset$-interpretable language) endowed with their full structure;
- the root is $\mathcal{M}$;
- any node admits relative resplendent quantifier elimination (in some $\emptyset$-interpretable language) to the set of its children;
- any two sorts in two different branches are orthogonal.

The idea is that one might be able to reduce certain questions on the structure $\mathcal{M}$ to the set of its leaves. Every node describes then an intermediate step. Reduction to a node would also have the advantage of being generalised to any enrichment of structure below the node.

In this text, we compute the burden (Definition 1.12) of the following examples in terms of the burden of the leaves. In [38] we also characterise stable embeddedness of elementary pairs of models in terms of stable embeddedness of elementary pairs of structures in the leaves.
Example. 1. If $\mathcal{M}_{0}, \mathcal{M}_{1}$ are arbitrary structures, both the direct product $\mathcal{M}_{0} \times \mathcal{M}_{1}$ and the disjoint union $\mathcal{M}_{0} \cup \mathcal{M}_{1}$ reduce to $\mathcal{M}_{0}$ and $\mathcal{M}_{1}$ (Fact 1.33):


We can of course keep going:

2. Let $\mathcal{K}_{a c}=\{K, \Gamma, k$, ac $: K \rightarrow k\}$ be a Henselian valued field of equicharacteristic 0 of valued group $\Gamma$, residue field $k$, and angular component ac. It admits the following reduction diagrams (Theorem of Pas):

3. Let $\mathcal{M}=\{A, B, C, \iota, \nu\}$ be a short exact sequence of abelian groups

$$
0 \longrightarrow A \xrightarrow{\iota} B \xrightarrow{\nu} C \longrightarrow 0 \text {, }
$$

seen as a three-sorted structure. Assume $A$ is a pure subgroup of $B$. It admits the following reduction diagram (Fact 1.74):


To get relative quantifier elimination, one has to consider interpretable maps from $B$ to $A / n A$, $n \geq 0$. The sort $A / n A$ are understood to be part of the induced structure on $A$.
4. Let $\mathcal{K}=\{K, \Gamma, k, \operatorname{RV}(K)\}$ be a Henselian valued field of equicharacteristic 0 , value group $\Gamma$, residue field $k$ and RV -sort $\mathrm{RV}(K)$ (definition in Paragraph 1.2.1). It admits the following reduction diagram (Fact 1.63 and Fact 1.74):

5. If $\mathcal{K}=\{K, \Gamma, k\}$ is a Henselian valued field of equicharacteristic 0 , where moreover the residue field $k$ is endowed with a structure $\left(k, \Gamma^{\prime}, k^{\prime}\right)$ of Henselian valued field of equicaracteristic 0 , and $\Gamma$ is endowed with a predicate for a convex subgroup $\Delta$. Then by Corollary 1.75 (and resplendence), we have the following reduction diagram:


## 1 Preliminaries

### 1.1 On pure model theory

We will assume the reader to be familiar with basic model theory concepts, and in particular with standard notations. One can refer to [37]. Symbols $x, y, z, \ldots$ will usually refer to tuples of variables, $a, b, c, \ldots$ to parameters. Capital letters $K, L, M, N, \ldots$ will refer to sets, and calligraphic letters $\mathcal{K}, \mathcal{L}, \mathcal{M}, \mathcal{N}, \ldots$ will refer to structures with respective base sets $K, L, M, N, \ldots$. If there is no ambiguity, we may respectively name a very saturated elementary extension with blackboard bold letters $\mathbb{K}, \mathbb{L}, \mathbb{M}, \ldots$ Languages will be denoted with a roman character $\mathrm{L}, \mathrm{L}^{\prime}, \mathrm{L}_{\text {Rings }}, \mathrm{L}_{\Gamma, k}$ etc.

In this section, we will consider any (possibly multi-sorted) first order language L , and an arbitrary L-structure $\mathcal{M}$.

### 1.1.1 Relative quantifier elimination and resplendence

We will use freely Rideau-Kikuchi's terminology about enrichment that we briefly recall now. The reader can refer to [30, Annexe A] for a more detailed exposition. This will allow us to give effortless generalisations to richer structures, or to simplify the notation, by reducing the language to the strict necessity for producing transfer principles.

First, let us recall two notions of relative quantifier elimination.
Definition 1.1. Let $\mathcal{M}$ be a multisorted structure in a language $L$, and consider $\Pi \cup \Sigma$ a partition of the set of sorts. We denote by $\left.\mathrm{L}\right|_{\Sigma}$ the language of all function symbols and relation symbols in L involving only sorts in $\Sigma$. Then, we say that

- $\mathcal{M}$ eliminates $\Pi$-quantifiers if every formula $\phi(x)$ is equivalent to a formula without quantifier in a sort in $\Pi$.
- $\mathcal{M}$ eliminates quantifiers relatively to $\Sigma$ if the theory of $\mathcal{M}^{\Sigma-\text { Mor }}$ - obtained by naming all $\left.\mathrm{L}\right|_{\Sigma}$-definable sets (without parameters) with a new predicate- eliminates quantifiers.

As observed in [30, Annexe A], $M$ eliminates quantifiers relatively to $\Sigma$, then it eliminates $\Pi$ quantifiers.

Definition 1.2. Let $\mathcal{M}$ be a multi-sorted structure in a language L , and let $\Sigma$ be a set of sorts in L .

- a language $\mathrm{L}_{e}$ containing L is said to be a $\Sigma$-enrichment of L if all new function symbols and relation symbols only involve the sorts in $\Sigma$ and the new sorts $\Sigma_{e}$ in $\mathrm{L}_{e} \backslash \mathrm{~L}$. An expansion $\mathcal{M}_{e}$ of $\mathcal{M}$ to $\mathrm{L}_{e}$ is called a $\Sigma$-enrichment of $\mathcal{M}$.
- $\Sigma$ is said to be closed if any relation symbol involving a sort in $\Sigma$ or any function symbol with a domain involving a sort in $\Sigma$ only involves sorts in $\Sigma$.

Fact 1.3 ([30]). Let $\mathcal{M}$ be a multisorted structure, and consider $\Pi \cup \Sigma$ a partition of the set of sorts. If $\Sigma$ is a closed set of sorts, then $\mathcal{M}$ eliminates $\Pi$-quantifiers if and only if $\mathcal{M}$ eliminates quantifiers relatively to $\Sigma$.

In the context of this text, these two notions of quantifier elimination will be often equivalent. Another consequence of closedness is the automatic resplendence of relative quantifier elimination:

Definition 1.4. Let $\mathcal{M}$ be a multi-sorted structure in a language L , and let $\Sigma$ be a set of sorts in L . We say that $\mathcal{M}$ eliminates quantifiers resplendently relatively to $\Sigma$ if for any $\Sigma$-enrichment $\mathcal{M}_{e}$ of $\mathcal{M}$, $\operatorname{Th}\left(\mathcal{M}_{e}\right)$ eliminates quantifiers relatively to $\Sigma \cup \Sigma_{e}$ (where $\Sigma_{e}$ is the set of new sorts in $\mathrm{M}_{e}$ ).

Fact 1.5 ([30, Proposition A.9]). Let $\mathcal{M}$ be a multi-sorted structure in a language L , and assume that $\operatorname{Th}(\mathcal{M})$ eliminates quantifiers relative to a closed set of sorts $\Sigma$. Then $\operatorname{Th}(\mathcal{M})$ eliminates quantifiers resplendently relatively to $\Sigma$.

Notice however that closedness does not characterise resplendence of relative quantifier elimination, as we will see later with pure short exact sequence of abelian groups. Let us introduce the notion of stable embedded definable sets and of pure sorts.

Definition 1.6. - A definable subset $D$ of $\mathcal{M}$ is called stably embedded if all definable subsets of $D^{n}, n \in \mathbb{N}$ can be defined with parameters in $D$.

- Two definable subsets $D$ and $D^{\prime}$ of $\mathcal{M}$ are called orthogonal if for all formulas

$$
\phi\left(x_{0}, \ldots, x_{n-1} ; x_{0}^{\prime}, \ldots, x_{n-1}^{\prime}, a\right)
$$

with parameters $a$ in $\mathcal{M}$, there are finitely many formulas $\theta_{i}\left(x_{0}, \ldots, x_{n-1}, a_{i}\right)$ and $\theta_{i}^{\prime}\left(x_{0}^{\prime}, \ldots, x_{n-1}^{\prime}, a_{i}^{\prime}\right)$, with $i<k$ and parameters $a_{0}, \ldots, a_{n-1}, a_{0}^{\prime}, \ldots, a_{n-1}^{\prime}$ in $\mathcal{M}$, such that

$$
\phi\left(D^{n}, D^{\prime n}, a\right)=\cup_{i<k} \theta_{i}\left(D^{n}, a_{i}\right) \times \theta_{i}^{\prime}\left(D^{\prime n}, a_{i}\right) .
$$

If $S$ is a sort, we use the following terminology in order to say that definable sets in $S$ can be given by formulas with parameters in $S$ and function/predicate symbols contained in $S$.

Definition 1.7. Consider a sort $S$ in an L-structure $\mathcal{M}$. We denote by $\mathrm{L} \mid S$ the language L restricted to function/predicate symbols which only involve $S$. We say that $S$ is pure or unenriched if definable subsets of $S$ (with parameters) are given by $\mathrm{L}_{S}(S)$-formulas - this is to say by a formula with parameters in $S$ and with function/predicate symbols only involving the sort $S$.

A pure sort $S$ can be seen as an $\mathrm{L}_{S_{S}}$-structure on its own. In particular, it is stably embedded. Purity of a sort $S$ is usually a simple corollary of quantifier elimination relative to $S$ and closedness of $S$ (see Fact 1.9).

Remark that the notion of closedness is syntactic, which is not ideal. One may indeed use another bi-interpretable language, where the sort is no longer closed, but where resplendent relative quantifier elimination still holds ${ }^{2}$. Here is a sightly improved version of purity which can replace the notion of closedness. It is, in a certain sense, less dependent of the language.

Definition 1.8. Consider $\mathcal{M}$ a structure. An imaginary sort $\mathcal{S}=(S, \ldots)$ endowed with an interpretable structure in a language $\mathrm{L}_{S}$ is called pure with control of parameters if every formula $\phi_{i}\left(x_{S}, b\right)$ where $x_{S}$ is a tuple of $S$-variables and $b$ is a tuple of parameters in M, is equivalent to a formula $\phi_{S}\left(x_{S}, \pi_{S}(t(b))\right)$ where $\pi_{S}$ is the canonical projection onto $S, \phi_{S}$ is an $\mathrm{L}_{S}$-formula and $t(x)$ is a tuple of L-terms.

The following is immediate:
Fact 1.9. Consider $\mathcal{M}$ an L -structure. If $S$ is a closed sort and $\mathcal{M}$ eliminates quantifier relative to $S$, then $S$ endowed with its induced structure in $\mathrm{L}_{\mid S}$ is pure with control of parameters. In particular, it is pure and stably embedded.

Proposition 1.10. Let $\mathcal{M}$ be an L-structure, and $\mathcal{S}$ an imaginary sort of arity $n$ with some interpretable structure given in a language $\mathrm{L}_{S}$. Assume that $\mathcal{M}$ has quantifier elimination and that $\mathcal{S}$ is a pure imaginary structure with control of parameters. The (multisorted) structure $\left\{\mathcal{M}, \mathcal{S}, \pi_{S}: M^{n} \rightarrow S\right\}$ in the language $\mathrm{L}_{2}:=\mathrm{L} \cup \mathrm{L}_{S} \cup\left\{\pi_{S}: M^{n} \rightarrow S\right\}$ admits quantifier elimination relatively to $\mathcal{S}$.

[^2]Notice that by definition, we have $\left.\mathrm{L}_{2}\right|_{S}=\mathrm{L}_{S}$. As $\mathcal{S}$ is by definition a closed sort in the language $\mathrm{L}_{2}$, this is in fact a characterisation of purity with control of parameters:

Corollary 1.11. An interpretable structure $\mathcal{S}$ is pure with control of parameters if and only if $\left\{\mathcal{M}, \mathcal{S}, \pi_{S}: M^{n} \rightarrow S\right\}$ admits quantifier elimination relative to $\mathcal{S}$.

Proof. We use the usual back-and-forth argument. Let $\mathcal{N}$ be a $|M|$-saturated model of the theory of $\mathcal{M}$ in the language $\mathrm{L}_{2}$. Let $f:\left(A, S_{A}\right) \rightarrow\left(B, S_{B}\right)$ be an isomorphism between a substructure $\left(A, S_{A}\right)$ of $\mathcal{M}$ and a substructure $\left(B, S_{B}\right)$ of $\mathcal{N}$. Assume that the restriction $\left.f\right|_{S}$ to $\mathcal{S}$ is elementary. We want to extend $f$ to an embedding of $\mathcal{M}$ into $\mathcal{N}$.
Step 0: We may assume that $S_{A}=S_{M}$.
Indeed, by elementarity of $\left.f\right|_{S}$, there exists an isomophism $\left.\tilde{f}\right|_{S}:\left.S_{M} \rightarrow \tilde{f}\right|_{S}\left(S_{M}\right) \subset S_{N}$ extending $\left.f\right|_{S}$. The union $\left.f \cup \tilde{f}\right|_{S}$ is a partial isomorphism as the sort $S$ is closed. Indeed, every quantifier-free formula $\phi(a, s)$ with parameters in $\left(A, S_{M}\right)$ can be written of the form:

$$
\bigvee \phi_{\mathrm{L}}(a) \wedge \phi_{\mathcal{S}}\left(s, \pi_{S}(t(a))\right)
$$

where $\phi_{\mathrm{L}}$ is an L-formula, $\phi_{\mathcal{S}}$ is an $\mathrm{L}_{S}$-formula and $t$ is a tuples of L-terms. As $A$ is a structure, all terms $t(a)$ are elements of $A$. It follows that $\left.f \cup \tilde{f}\right|_{S}$ preserves these formulas.
Step 1: We may assume that $A=M$ and thus conclude the proof.
Indeed, let $a \in M \backslash A$. We denote by $p(x)$ the quantifier free type of $a$ over $A$. We want an appropriate answer for $a$, i.e. an element $\tilde{f}(a)$ of $\mathcal{N}$ satisfying the set of formulas:

$$
\left\{\phi(x, f(b), f(s)) \mid \phi(x, b, s) \in p(x), b \in A, s \in S_{M}\right\}
$$

By compactness, it suffices to show that it is finitely consistent. Consider a formula

$$
\phi(x, b, s) \in p(x)
$$

where $b \in A$ and $s \in S_{M}$. As $S$ is pure with control of parameters, the formula

$$
\exists x \phi\left(x, b, y_{S}\right)
$$

is equivalent to an $\mathrm{L}_{S}\left(S_{M}\right)$-formula $\psi_{S}\left(t(b), y_{S}\right.$ ) (with a tuple of L-terms $t(y)$ ).
The formula

$$
\theta(y)=\forall y_{S} \psi_{S}\left(t(y), y_{S}\right) \Leftrightarrow \exists x \phi\left(x, y, y_{S}\right)
$$

is interpreted in the language L by a formula $\Sigma(y)$. As $\mathcal{M}$ has quantifier elimination in the language L, we may assume that $\Sigma(y)$ is quantifier-free. We have:

$$
\begin{aligned}
\mathcal{M} & =\Sigma(b) \\
\mathcal{M} & =\psi_{S}(t(b), s)
\end{aligned}
$$

As $f$ respects quantifier free-formula and $\left.f\right|_{S_{M}}$ respects $L_{S}\left(S_{M}\right)$-formula, we have

$$
\begin{array}{r}
\mathcal{N} \models \Sigma(f(b)), \\
\mathcal{N} \vDash \psi_{S}(f(t(b)), f(s)) .
\end{array}
$$

Of course, $f(t(b))=t(f(b))$. We get: $\mathcal{N} \vDash \exists x \phi(x, f(b), f(s))$. This concludes our proof.
As an example, we treat the question of quantifier elimination in the field of $p$-adics in a two-sorted language of valued fields. This is a well known result, but we are not aware of a reference.

Example. Consider the theory $T$ of the $p$-adics $\mathbb{Q}_{p}$ for some $p$. By Macintyre's theorem [26], it admits quantifier elimination in the language $\mathrm{L}_{M a c}:=\mathrm{L}_{\text {rings }} \cup\left\{P_{n}\right\}_{n<\omega}$ where the predicate $P_{n}$ interprets the $n^{\text {th }}$-powers.

- The value group $\Gamma$, simply considered as a set, is not pure. Indeed, the theory $T$ in the language $\mathrm{L}_{M a c} \cup\{\Gamma\} \cup\{\operatorname{val}\}$ (no structure on $\Gamma$ ) does not eliminate the quantifiers in the formula encoding the addition:

$$
\phi\left(x_{\Gamma}, y_{\Gamma}, z_{\Gamma}\right) \equiv \exists x, y \in K \operatorname{val}(x)=x_{\Gamma} \wedge \operatorname{val}(y)=y_{\Gamma} \wedge \operatorname{val}(x y)=z_{\Gamma}
$$

where $x_{\Gamma}, y_{\Gamma}, z_{\Gamma}$ are variables in $\Gamma$.

- By Bélair's theorem $\left[6\right.$, Theorem 5.1], the structure $\left\{\mathbb{Q}_{p}, \mathcal{O}_{n}, \Gamma\right.$, val : $\left.\mathbb{Q}_{p} \rightarrow \mathbb{Z}, \mathrm{ac}_{n}: \mathbb{Q}_{p} \rightarrow \mathcal{O}_{n}\right\}$ enriched with angular components (see Paragraph 1.2.1) eliminates quantifiers in the sort for $\mathbb{Q}_{p}$. It results that the value group $(\Gamma,+, 0,<, \infty)$-as an imaginary sort of $\mathbb{Q}_{p}$ in the ring languageis pure with control of parameters. Then, by Corollary $1.11, T$ eliminates quantifiers relative to $\Gamma$ in the language

$$
\mathrm{L}_{2}:=\mathrm{L}_{M a c} \cup\{\Gamma,<,+, 0, \infty\} \cup\{\operatorname{val}\}
$$

- To get full elimination of quantifiers, one only needs to eliminate quantifiers in $\{\Gamma,<,+, 0, \infty\}$. So the theory $T$ eliminates quantifiers in the language $\mathrm{L}_{M a c} \cup\left\{\Gamma,<,+, P_{\Gamma, n}, 0,1, \infty\right\} \cup\{\operatorname{val}\}$ where $P_{\Gamma, n}$ interprets the set of values divisible by $n$.


### 1.1.2 Burden of a theory

In [32], Shelah defined the notion of burden as an invariant cardinal $\kappa_{i n p}$ and implicitly defined the tree property of the second kind. A theory which does not satisfy it is called $\mathrm{NTP}_{2}$. Interest in the class of $\mathrm{NTP}_{2}$ theories grew after the success of stability theory and with the necessity of extending methods to unstable contexts. In [10], Chernikov and Kaplan studied the forking relation in $\mathrm{NTP}_{2}$ theories, establishing notably that types over models fork if and only if they divide. In [8], Chernikov continued the study of $\mathrm{NTP}_{2}$ theories, establishing in particular a criterion with indiscernible sequences and the sub-multiplicativity of the burden.

We recall here a definition of burden, some of the results cited above and give some important lemmas required for the proof of Theorem 3.12. We will give a second definition ( slightly different) of the burden in order to formalise a convention due to Adler [1].

Definition 1.12. Let $\lambda$ be a cardinal. For all $i<\lambda, \phi_{i}\left(x, y_{i}\right)$ is L-formula where $x$ is a common tuple of free variables, $b_{i, j}$ are elements of $\mathbb{M}$ of size $\left|y_{i}\right|$ and $k_{i}$ is a positive natural number. Finally, let $p(x)$ be a partial type. We say that $\left\{\phi_{i}\left(x, y_{i}\right),\left(b_{i, j}\right)_{j \in \omega}, k_{i}\right\}_{i<\lambda}$ is an inp-pattern of depth $\lambda$ in $p(x)$ if:

1. for all $i<\lambda$, the $i^{\text {th }}$ row is $k_{i}$-inconsistent: any conjunction $\bigwedge_{l=1}^{k_{i}} \phi_{i}\left(x, b_{i, j_{l}}\right)$ with $j_{1}<\cdots<j_{k_{i}}<\omega$, is inconsistent.
2. all (vertical) paths are consistent: for every $f: \lambda \rightarrow \omega$, the set $\left\{\phi_{i}\left(x, b_{i, f(i)}\right)\right\}_{i<\lambda} \cup p(x)$ is consistent.

Most of the time, we will not mention the $k_{i}$ 's and only say that the rows are finitely inconsistent.
Definition 1.13. - Let $p(x)$ be a partial type. The burden of $p(x)$, denoted by $\operatorname{bdn}(p(x))$, is the cardinal defined as the supremum of the depths of inp-patterns in $p(x)$. If $C$ is a small set of parameters, we write $\operatorname{bdn}(a / C)$ instead of $\operatorname{bdn}(\operatorname{tp}(a / C))$.

- The cardinal $\sup _{S \in \mathcal{S}} \operatorname{bdn}\left(\left\{x_{S}=x_{S}\right\}\right)$ where $x_{S}$ is a single variable from the sort $S$ and $\mathcal{S}$ is the set of sorts, is called the burden of the theory $T$, and it is denoted by $\kappa_{\text {inp }}^{1}(T)$ or by $\operatorname{bdn}(T)$. The theory $T$ is said to be inp-minimal if $\kappa_{\text {inp }}^{1}(T)=1$.
- More generally, for $\lambda$ a cardinal, we denote by $\kappa_{\text {inp }}^{\lambda}(T)$ the supremum of $\operatorname{bdn}(\{x=x\})$ where $|x|=\lambda$ and variables run in all sorts $S \in \mathcal{S}$. We always have $\kappa_{\text {inp }}^{\lambda}(T) \geq \lambda \cdot \kappa_{\text {inp }}^{1}(T)$. In particular, if models of $T$ are infinite, $\kappa_{\text {inp }}^{\lambda}(T) \geq \lambda$.
- A formula $\phi(x, y)$ has $\mathrm{TP}_{2}$ if there is an inp-pattern of the form $\left\{\phi(x, y),\left(b_{i, j}\right)_{j<\omega}, k_{i}\right\}_{i<\omega}$. Otherwise, we say that $\phi(x, y)$ is $\mathrm{NTP}_{2}$.
- The theory $T$ is said $\mathrm{NTP}_{2}$ if $\kappa_{\text {inp }}^{1}(T)<\infty$. Equivalently, $T$ is $\mathrm{NTP}_{2}$ if and only if there is no $\mathrm{TP}_{2}$ formula. (See [8, Remark 3.3])

In all the notation above, we may replace $T$ by $\mathcal{M}$. In [8], Chernikov proves the following:
Fact 1.14 (Sub-multiplicativity). Let $a_{1}, a_{2} \in \mathbb{M}$. If there is an inp-pattern of depth $\kappa_{1} \times \kappa_{2}$ in $\operatorname{tp}\left(a_{1} a_{2} / C\right)$, then either there is an inp-pattern of depth $\kappa_{1}$ in $\operatorname{tp}\left(a_{1} / C\right)$ or there is an inp-pattern of depth $\kappa_{2}$ in $\operatorname{tp}\left(a_{2} / a_{1} C\right)$.

As a corollary, for $n<\omega$, we have $\kappa_{i n p}^{n}(T)+1 \leq\left(\kappa_{i n p}^{1}(T)+1\right)^{n}$ and then $\kappa_{i n p}^{n}(T)=\kappa_{i n p}^{1}(T)=\operatorname{bdn}(T)$ as soon as one of these cardinals is infinite.

If the reader knows the notion of dp-rank of a theory $T$, usually denoted by dp-rank $(T)$, let us say the following: it admits as well a similar definition in term of depth of ict-patterns and it has been showed that a theory $T$ is NIP if and only if the depth of $i c t$-pattern is bounded by some cardinal (see [1] for more details). In this paper, the reader only needs to know that the notions of dp-rank and burden coincide in NIP theories. If they are not familiar with the notion of dp-rank, they may take it as a definition.

Fact $1.15([1$, Proposition 10]). Let $T$ be an NIP theory, and $p(x)$ a partial type. Then $d p-r a n k(p(x))=$ $\operatorname{bdn}(p(x))$.

The previous fact is only stated with partial type $p(x)=\{x=x\}$, but the proof is the same.
Example. Any quasi-o-minimal theory is inp-minimal (see e.g. [35, Theorem A.16]). In particular, $\{\mathbb{Z}, 0,+,<\}$ is inp-minimal.

We give below, for any cardinal $\lambda$, an example of a structure of burden $\lambda$. First, we want to give some tools in order to 'manipulate' inp-patterns.

The definition of burden of a theory, as many other notion of complexity, gives to unary sets an important role. But one has to notice that the notion of unary set is syntactic, and is not preserved under bi-interpretability:

Remark 1.16. For $n \in \mathbb{N}$, one can consider the multisorted structure $\left(\mathcal{M}^{n}, \mathcal{M}, p_{i}, i<n\right)$ where $p_{i}: \mathcal{M}^{n} \rightarrow \mathcal{M},\left(a_{0}, \ldots, a_{n-1}\right) \mapsto a_{i}$ is the projection to the $i^{\text {th }}$ coordinate. If we denote its theory by $T^{n}$, then we clearly have $\kappa_{\text {inp }}^{n}(T)=\kappa_{\text {inp }}^{1}\left(T^{n}\right)$.

To clarify, let us introduce the following terminology:
Definition 1.17. Let $\mathcal{M}$ and $\mathcal{N}$ be two structures. We say that $\mathcal{N}$ is interpretable on a unary set in $\mathcal{M}$ if there is a bijection $f: N \rightarrow D / \sim$ where $D$ is a unary definable set in $\mathcal{M}, \sim$ is a definable equivalence relation, and the pull-back in $\mathcal{M}$ of any graph of function and relation of $\mathcal{N}$ is definable. The structures $\mathcal{M}$ and $\mathcal{N}$ will be said bi-interpretable on unary sets ${ }^{3}$ if $\mathcal{N}$ is interpretable on a unary set in $\mathcal{M}$ and $\mathcal{M}$ is interpretable on a unary set in $\mathcal{N}$.

[^3]We will work up to bi-interpretability on unary sets, meaning in particular that the main results of this text will only depend on the structure that we want to consider and not on the language.

Fact 1.18. Let $\mathcal{M}$ and $\mathcal{N}$ be two structures, and assume that $\mathcal{N}$ is interpretable on a unary set in $\mathcal{M}$, then $\operatorname{bdn}(\mathcal{N}) \leq \operatorname{bdn}(\mathcal{M})$. In particular, if $\mathcal{M}$ and $\mathcal{N}$ are bi-interpretable on unary sets, then $\operatorname{bdn}(\mathcal{M})=\operatorname{bdn}(\mathcal{N})$.

For example, $\{\mathbb{Z}, 0,+,<\}$ does not interpret

$$
\left\{\mathbb{Z} \times \mathbb{Z},(\mathbb{Z}, 0,+,<), \pi_{1}: \mathbb{Z} \times \mathbb{Z} \rightarrow \mathbb{Z}, \pi_{2}: \mathbb{Z} \times \mathbb{Z} \rightarrow \mathbb{Z}\right\}
$$

on a unary set (the first being of burden 1 , the second being of burden at least 2). However, if $k$ is an imperfect field, we will see that $k$ interprets $\left\{k \times k,(k, 0,1,+, \cdot), \pi_{1}: k \times k \rightarrow k, \pi_{2}: k \times k \rightarrow k\right\}$ on a unary set.

Lemmas on inp-patterns Let L be any first order language, $\mathcal{M}$ a L-structure of base set $M$ and let $\lambda$ be a cardinal.

Definition 1.19. - A sequence $\left(b_{j}\right)_{j \in \lambda}$ of (tuples of) elements of $M$ is indiscernible over a subset $A \subset M$ if for every $n \in \mathbb{N}$ and every formula $\phi\left(x_{0}, \ldots, x_{n-1}, a\right)$ with parameters $a \in A$, we have

$$
\mathcal{M} \vDash \phi\left(b_{0}, \ldots, b_{n-1}, a\right) \Leftrightarrow \phi\left(b_{j_{0}}, \ldots, b_{j_{n-1}}, a\right)
$$

for every $j_{0}<\cdots<j_{n-1} \in \lambda$.

- An array $\left(b_{i, j}\right)_{i \in \lambda, j \in \omega}$ is mutually indiscernible if every line $\left(b_{i, j}\right)_{j \in \omega}$ is indiscernible over $\left\{b_{k, j}\right\}_{k \neq i, k \in \lambda, j<\omega}$.
We will intensively use the following fact:
Fact 1.20 ([8, Lemma 2.2]). If $p(x)$ is a partial type and if $\left\{\phi_{i}\left(x, y_{i}\right),\left(b_{i, j}\right)_{j \in \omega}, k_{i}\right\}_{i<\lambda}$ is an inp-pattern in $p(x)$, there is an inp-pattern $\left\{\phi_{i}\left(x, y_{i}\right),\left(\tilde{b}_{i, j}\right)_{j \in \omega}, k_{i}\right\}_{i<\lambda}$ in $p(x)$ with a mutually indiscernible array $\left(\tilde{b}_{i, j}\right)_{i<\lambda, j<\omega}$.

We will now present some easy lemmas, which we will later use. They give us tools to 'transform' inp-patterns into simpler ones which are easier to analyse.

Lemma 1.21. Let $\left\{\phi_{i}\left(x, y_{i}\right),\left(a_{i, j}\right)_{j<\omega}, k_{i}\right\}_{i<\lambda}$ be an inp-pattern with $\left(a_{i, j}\right)_{i<\lambda, j<\omega}$ mutually indiscernible. Assume for every $i<\lambda, \phi_{i}\left(x, a_{i, 0}\right)$ is equivalent to some formula $\psi_{i}\left(x, b_{i, 0}\right)$ with parameter $b_{i, 0}$. Then we may extend $\left(b_{i, 0}\right)_{i<\lambda}$ to an mutually indiscernible array $\left(b_{i, j}\right)_{i<\lambda, j<\omega}$ such that

$$
\left\{\psi_{i}\left(x, y_{i}\right),\left(b_{i, j}\right)_{j<\omega}, k_{i}\right\}_{i<\lambda}
$$

is an inp-pattern.
Proof. By 1-indiscernibility, we find $b_{i, j}$ such that $\phi_{i}\left(\mathbb{M}, a_{i, j}\right)=\psi_{i}\left(\mathbb{M}, b_{i, j}\right)$. Then, the statement is clear.

Remark 1.22. Let $D$ be a stably embedded definable set in $\mathbb{M}$, and $\left\{\phi_{i}\left(x, y_{i, j}\right),\left(a_{i, j}\right)_{j<\omega}, k_{i}\right\}_{i<\lambda}$ an inppattern in $D$. This in particular implies that solutions of paths can be found in $D$ but the parameters $\left(a_{i, j}\right)$ may not belong to $D$. Using the previous lemma, we may actually assume that this is the case. It follows that $D$ endowed with the induced structure is at least of burden $\lambda$.

The next lemma shows that one can 'eliminate' disjunction symbols in inp-patterns. A direct consequence is that if the theory has quantifier elimination, then we may assume that formulas of inp-patterns are conjunctions of atomic and negation of atomic formulas.

Lemma 1.23. Let $\left\{\phi_{i}\left(x, y_{i, j}\right),\left(a_{i, j}\right)_{j<\omega}, k_{i}\right\}_{i<\lambda}$ be an inp-pattern with $\left(a_{i, j}\right)_{i<\lambda, j<\omega}$ mutually indiscernible. Assume that $\phi_{i}\left(x, y_{i, j}\right)=\bigvee_{l \leq n_{i}} \psi_{l, i}\left(x, y_{i, j}\right)$. Then there exists a sequence of natural numbers $\left(l_{i}\right)_{i<\lambda}$ such that $l_{i} \leq n_{i}$ and

$$
\left\{\psi_{l_{i}, i}\left(x, y_{i, j}\right),\left(a_{i, j}\right)_{j<\omega}, k_{i}\right\}_{i<\lambda}
$$

is an inp-pattern.
Proof. Let $d \models\left\{\phi_{i}\left(x, a_{i, 0}\right)\right\}_{j<\lambda}$. For every $i<\lambda$, let $l_{i} \leq n_{i}$ be such that $d \models \psi_{l_{i}, i}\left(x, a_{i, 0}\right)$. By the mutual-indiscernibility of $\left(a_{i, j}\right)_{i<\lambda, j<\omega}$, every path of the pattern $\left\{\psi_{l_{i}, i}\left(x, y_{i, j}\right),\left(a_{i, j}\right)_{j<\omega}, k_{i}\right\}_{i<\lambda}$ is consistent. The inconsistency of the rows follows immediately from the inconsistency of the rows of the initial pattern.

This lemma is particularly useful in order to compute the burden, as shown in the following example.
Example. - Let $\mathrm{L}=\{R, B\}$ be the language with two binary predicates, and let $\mathcal{M}$ be a set with two cross-cutting equivalence relations with infinitely many infinite classes (for all $a$ and $b$, there are infinitely many $c$ such that $R(a, c)$ and $B(b, c))$.


Then we have $\operatorname{bdn}(\mathcal{M})=2$.

- More generally, given a cardinal $\lambda$, we can consider a set $\mathcal{M}_{\lambda}$ equipped with $\lambda$-many cross cutting equivalence relations (in a language with $\lambda$-many binary predicates). Then, $\operatorname{bdn}\left(\mathcal{M}_{\lambda}\right)$ is of burden exactly $\lambda$.

Proof. For simplicity, we only show the first case (when $\lambda=2$ ). The general case can be prove the same way. We also leave to the reader the proof that $(\mathcal{M}, B, R)$ eliminates quantifiers and show that $\operatorname{bdn}(\mathcal{M})=2$. First, let us give an inp-pattern of depth 2. Consider $\left(a_{0, j}\right)_{j<\omega}$ in distinct $R$-equivalence classes and $\left(a_{1, j}\right)_{j<\omega}$ in distinct $B$-equivalence classes. Then one sees that

$$
\left\{R\left(x, y_{0}\right), B\left(x, y_{1}\right),\left(a_{i, j}\right)_{j<\omega}, 2\right\}_{i<2}
$$

is an inp-pattern of depth 2. Indeed:

- Rows are 2 -inconsistent: if $j, j^{\prime}<\omega$ are distinct, then $R\left(x, a_{0, j}\right) \wedge R\left(x, a_{0, j^{\prime}}\right)$ is inconsistent and $B\left(x, a_{1, j}\right)^{B}\left(x, a_{1, j^{\prime}}\right)$ is inconsistent.
- Paths are consistent: there are elements $c_{j_{0}, j_{1}}$ such that $R\left(a_{0, j_{0}}, c_{j_{0}, j_{1}}\right)$ and $B\left(a_{1, j_{1}}, c_{j_{0}, j_{1}}\right)$ for all $j_{0}, j_{1}<\omega$.
We show now that any inp-pattern of depth $>2$ must be of depth exactly 2 , and is - after some manipulation - of this form. Consider an inp-pattern

$$
\left\{\phi_{i}\left(x, y_{i}\right),\left(a_{i, j}\right)_{j<\omega}, k_{i}\right\}_{i<k}
$$

of depth $k>1$ with a mutual indiscernible array $\left(a_{i, j}\right)_{i<k, j<\omega}$ (we used Fact 1.20). By quantifier elimination and elimination of the disjunction (Lemma 1.23), we may assume that for all $i<k$, the
formula $\phi_{i}\left(x, y_{i}\right)$ is a conjunction of the atomic formulas $R\left(x, y_{i}^{l}\right), B\left(x, y_{i}^{l}\right), x=y_{i}^{l}\left(\right.$ for $\left.l<\left|y_{i}\right|\right)$ and their negation. In fact, we may assume that there is at most one positive occurrence of $B\left(x, y_{i}^{l}\right)$ in $\phi_{i}\left(x, y_{i}\right)$ and if one does occur, their is no other occurrence of $B$ in $\phi_{i}\left(x, y_{i}\right)$. Indeed, write for instance $\phi_{i}\left(x, y_{i}\right)=\phi_{i}^{\prime}\left(x, y_{i}\right) \wedge B\left(x, y_{i}^{0}\right) \wedge(\neg) B\left(x, y_{i}^{1}\right)$. Since the formula is consistent and $B$ is an equivalence relation, $\phi_{i}\left(x, a_{i, 0}\right)$ is equivalent to $\phi_{i}^{\prime}\left(x, a_{i, 0}\right) \wedge B\left(x, a_{i, 0}^{0}\right)$ and we may replace $\phi_{i}\left(x, y_{i}\right)$ by $\phi_{i}^{\prime}\left(x, y_{i}\right) \wedge B\left(x, y_{i}^{0}\right)($ Lemma 1.21).

Now, assume that there is a negative occurrence of $B$ in $\phi_{0}\left(x, y_{i}\right)$, say $\neg B\left(x, y_{0}^{0}\right)$. By the previous manipulation, their is no positive occurrence of $B$. Write $\phi_{0}\left(x, y_{0}\right)=\phi_{0}^{\prime}\left(x, y_{0}\right) \wedge \neg B\left(x, y_{0}^{0}\right)$. By writing that the line is $k_{0}$-inconsistent, we get $\bigwedge_{j<k_{0}} \phi_{0}^{\prime}\left(x, a_{0, j}\right) \vdash \bigvee_{j<k_{0}} B\left(x, a_{0, j}^{0}\right)$. If $\bigwedge_{j<k_{0}} \phi_{0}^{\prime}\left(x, a_{0, j}\right)$ is consistent, there are $0<j_{0}, j_{1}<k_{0}$ such that $\phi_{0}^{\prime}\left(x, a_{0, j_{0}}\right) \vdash B\left(x, a_{0, j_{1}}^{0}\right)$, which is a contradiction. Thus, $\left\{\phi_{0}^{\prime}\left(x, y_{0}\right),\left(a_{0, j}\right)_{j<\omega}\right\}$ is already $k_{0}$-inconsistent. We may remove all negative occurrence of $B$ in $\phi_{0}\left(x, y_{0}\right)$, as the line will still be $k_{0}$-inconsistent. Similarly, we remove all occurrences of $x \neq y_{0}^{l}$ and $\neg R\left(x, y_{i}^{l}\right)$ in $\phi_{0}\left(x, y_{0}\right)$ and repeat the same process for all formulas $\phi_{i}\left(x, y_{i}\right), 0<i<k$. It means that we may assume that $\phi_{i}\left(x, y_{i}\right)$ is one of the following four formulas: $R\left(x, y_{i}^{0}\right), B\left(x, y_{i}^{0}\right), B\left(x, y_{i}^{0}\right) \wedge R\left(x, y_{i}^{1}\right)$ or $B\left(x, y_{i}^{0}\right) \wedge R\left(x, y_{i}^{0}\right)$. Now, it is easy to see, using consistency of paths, that our pattern must be of depth 2 and of the form

$$
\left\{R\left(x, y_{0}^{0}\right), B\left(x, y_{1}^{0}\right),\left(a_{i, j}\right)_{j<\omega}\right\}
$$

'Elimination' of conjunction symbols may happen in more specific context. Notably:
Proposition 1.24. Let $\mathcal{K}$ and $\mathcal{H}$ be two structures, and consider the multisorted structure $\mathcal{G}$ :

$$
\mathcal{G}=\left\{K \times H, \mathcal{K}, \mathcal{H}, \pi_{K}: K \times H \rightarrow K, \pi_{H}: K \times H \rightarrow H\right\}
$$

called the direct product structure (where $\pi_{K}$ and $\pi_{H}$ are the natural projections). Then $\mathcal{G}$ eliminates quantifiers relative to $\mathcal{K}$ and $\mathcal{H}$, and $\mathcal{K}$ and $\mathcal{H}$ are orthogonal and stably embedded within $\mathcal{G}$.


We have

$$
\operatorname{bdn}(\mathcal{G})=\operatorname{bdn}(\mathcal{K})+\operatorname{bdn}(\mathcal{H})
$$

We prove an obvious generalisation for product of more than two structures in the next paragraph.
Proof. Relative quantifier elimination, stable embeddedness and orthogonality are rather obvious. The inequality $\operatorname{bdn}(\mathcal{G}) \geq \operatorname{bdn}(\mathcal{K})+\operatorname{bdn}(\mathcal{H})$ is easy but we give a detailed proof. Let $\left\{\phi_{i}\left(x_{K}, y_{i}\right),\left(a_{i, j}\right)_{j<\omega}\right\}_{i \in \lambda_{1}}$ be an inp-pattern in $\mathcal{K}$ and $\left\{\psi_{i}\left(x_{H}, y_{i}\right),\left(b_{i, j}\right)_{j<\omega}\right\}_{i \in \lambda_{2}}$ an inp-pattern in $\mathcal{H}$. Then

$$
\left\{\phi_{i}\left(\pi_{K}\left(x_{K}, x_{H}\right), y_{i}\right),\left(a_{i, j}\right)_{j<\omega}\right\}_{i \in \lambda_{1}} \cup\left\{\psi_{i}\left(\pi_{H}\left(x_{K}, x_{H}\right), y_{i}\right),\left(b_{i, j}\right)_{j<\omega}\right\}_{i \in \lambda_{2}}
$$

is an inp-pattern in $\mathcal{G}$ of depth $\lambda_{1}+\lambda_{2}$. Indeed, first notice that inconsistency of each rows is clear. Secondly, take a path $f: \lambda_{1} \sqcup \lambda_{2} \rightarrow \omega$. There is an element $d_{K} \in K$ satisfying $\left\{\phi_{i}\left(x_{K}, a_{i, f(i)}\right)\right\}_{i \in \lambda_{1}}$ and an element $d_{H} \in H$ satisfying $\left\{\psi_{i}\left(x_{H}, b_{i, f(i)}\right)\right\}_{i \in \lambda_{2}}$. Then, the element $d=\left(d_{K}, d_{H}\right)$ of $\mathcal{G}$ is a solution of the pattern along the path $f$.

For the other inequality, let $\left\{\theta_{i}\left(x, y_{i, j}\right),\left(c_{i, j}\right)_{j<\omega}, k_{i}\right\}_{i<\lambda}$ be an inp-pattern in $\mathcal{G}$, with $\left(c_{i, j}\right)_{i<\lambda, j<\omega}$ mutually indiscernible. We may assume $\theta_{i}\left(x, c_{i, j}\right)$ is of the form $\phi_{i}\left(x_{K}, a_{i, j}\right) \wedge \psi_{i}\left(x_{H}, b_{i, j}\right)$ where $x_{K}=$
$\pi_{K}(x), x_{H}=\pi\left(x_{H}\right), c_{i, j}=a_{i, j}{ }^{\wedge} b_{i, j}, \phi_{i}\left(x_{K}, a_{i, j}\right)$ is a $\mathcal{K}$-formula and $\psi_{i}\left(x_{H}, b_{i, j}\right)$ is a $\mathcal{H}$-formula. Indeed, let $d \models\left\{\theta_{i}\left(x, c_{i, 0}\right)\right\}_{j<\lambda}$, by orthogonality, $\theta_{i}\left(x, c_{i, 0}\right)$ is equivalent to a formula of the form:

$$
\bigvee_{k<n_{i}} \phi_{i, k}\left(x_{K}, a_{i, 0}\right) \wedge \psi_{i, k}\left(x_{H}, b_{i, 0}\right)
$$

Then we conclude by using Lemmas 1.21 and 1.23. For every $i$, at least one of the sets $\left\{\phi_{i}\left(x_{K}, a_{i, j}\right)\right\}_{j<\omega}$ and $\left\{\psi_{i}\left(x_{H}, b_{i, j}\right)\right\}_{j<\omega}$ is $k_{i}$-inconsistent (by indiscernibility of $\left.\left(c_{i, j}\right)_{j}\right)$. We may "eliminate" the conjunction as well and assume that every line is an $\mathrm{L}_{K}$-formula or an $\mathrm{L}_{H}$-formula. We conclude that $\lambda \leq \operatorname{bdn}(\mathcal{K})+\operatorname{bdn}(\mathcal{H})$.

Together with Fact 1.18, we get more generally:
Fact 1.25. Let $M=(A, C, \ldots)$ be a many-sorted structure. Assume that $A$ and $C$ are orthogonal and stably embedded in $M$. Then we have $\operatorname{bdn}(A \times C)=\operatorname{bdn}(A)+\operatorname{bdn}(C)$.

Let us finish this paragraph with one more lemma:
Lemma 1.26. Let $D$ and $D^{\prime}$ two type-definable sets respectively given by the partial types $p(x)$ and $p^{\prime}(x)$ and let $f: D \rightarrow D^{\prime}$ be a surjective finite to one type-definable function. Then we have $\operatorname{bdn}(D):=$ $\operatorname{bdn}(p(x))=\operatorname{bdn}\left(p^{\prime}(x)\right)=: \operatorname{bdn}\left(D^{\prime}\right)$.

Proof. We may assume that $D$ and $D^{\prime}$ are definable, the general case can be similarly deduced. Let $\left\{\phi_{i}^{\prime}\left(x^{\prime}, y_{i}\right),\left(a_{i, j}\right)_{j<\omega}, k_{i}\right\}_{i<\lambda}$ be an inp-pattern in $D^{\prime}$. Clearly, $\left\{\phi_{i}^{\prime}\left(f(x), y_{i}\right),\left(a_{i, j}\right)_{j<\omega}, k_{i}\right\}_{i<\lambda}$ is an inppattern in $D$. Hence $\operatorname{bdn}(D) \geq \lambda$. Conversely, let $\left\{\phi_{i}\left(x, y_{i}\right),\left(a_{i, j}\right)_{j<\omega}, k_{i}\right\}_{i<\lambda}$ be an inp-pattern of depth $\lambda$ in $D$. Consider the pattern

$$
\left\{\phi_{i}^{\prime}\left(x^{\prime}, y_{i, j}\right),\left(a_{i, j}\right)_{j<\omega}\right\}_{i<\lambda},
$$

where

$$
\phi_{i}^{\prime}\left(x^{\prime}, a_{i, j}\right) \equiv \exists x x \in D \wedge x^{\prime}=f(x) \wedge \phi\left(x, a_{i, j}\right)
$$

Clearly every path is consistent. Assume for some $i<\lambda$, the row $\left\{\phi_{i}^{\prime}\left(x^{\prime}, a_{i, j}\right)\right\}_{j<\omega}$ is consistent, witnessed by some $h^{\prime}$. Note that $h^{\prime}$ is in $D^{\prime}$. By the pigeonhole principle, there are $h \in D$ and an infinite subset $J$ of $\omega$ such that $f(h)=h^{\prime}$ and $h \models\left\{\phi_{i}\left(x, a_{i, j}\right)\right\}_{j \in J}$, contradiction. It follows that $\left\{\phi_{i}\left(x^{\prime}, y_{i, j}\right),\left(a_{i, j}\right)_{j<\omega}\right\}_{i<\lambda}$ is an inp-pattern in $D^{\prime}$. We conclude that $\operatorname{bdn}\left(D^{\prime}\right) \geq \lambda$.

### 1.1.3 More on burden and strength

We will formally introduce a well known convention with respect to the burden, which consists of writing $\operatorname{bdn}(\mathcal{M})=\lambda_{-}$for a limit cardinal $\lambda$ if $\mathcal{M}$ admits inp-patterns of depth $\mu$ for all $\mu<\lambda$, but no inp-pattern of depth $\lambda$. It has been introduced in [1], and has the advantage to emphasising a relevant distinction. If the reader is not interested by such subtleties, they may move to the next subsection. Proposition 1.34 might be interesting on its own, as it corresponds to the 'baby case' for the difficulty that we will encounter for mixed-characteristic Henselian valued fields. One can refer to [1] for this paragraph.

Definition 1.27. We define the ordered class $\left(\operatorname{Card}^{\star},<\right)$ as the linear order obtained from the ordered class of cardinals (Card, $<$ ) by adding for any limit cardinal $\lambda$ a new element $\lambda_{-}$(called 'lambda minus'). This new element comes immediately before $\lambda: \lambda_{-}<\lambda$ and if $\mu \in$ Card $^{\star}$ with $\mu<\lambda$, then $\mu \leq \lambda_{-}$. In addition to the natural injection Card $\hookrightarrow$ Card $^{\star}$, we define the actualisation map act : Card ${ }^{\star} \rightarrow$ Card as the map such that $\operatorname{act}\left(\lambda_{-}\right)=\lambda$ for every limit cardinal $\lambda$, and $\operatorname{act}(\kappa)=\kappa$ for any cardinal $\kappa \in$ Card. It will be convenient to also set $\kappa_{-}=\lambda$ when $\kappa=\lambda^{+} \in$ Card is a successor cardinal.

If $\lambda$ is a limit cardinal, one should think $\lambda$ as an 'actual' lambda and $\lambda_{-}$as a 'potential' lambda. We don't change our notion of cardinality of a set. As we will see, this definition of Card ${ }^{\star}$ is motivated by the burden, i.e. by a notion of dimension. It also motivates to (partially) extend the arithmetic operations of Card to Card*. We will have to answer any question of the form: should the cardinal $\aleph_{0} \cdot \aleph_{\omega-}$ be $\aleph_{\omega-}$ or $\aleph_{\omega}$ ? As the definitions themselves appear to be a bit technical, we prefer to first give intuition to the reader with a small digression on graphs.

Graphs and cliques We consider symmetric graphs in the language $\mathrm{L}=\{R\}$. We denote by $K_{\kappa}$ the complete graph on $\kappa$-many vertices, for $\kappa$ a cardinal in Card. Given a graph $\mathcal{G}$, we denote by $C(\mathcal{G})$ the cardinal in Card*:

$$
C(\mathcal{G})=\left\{\begin{array}{l}
\kappa \in \operatorname{Card} \text { if } K_{\kappa} \text { embeds in } \mathcal{G} \text { and } K_{\kappa^{+}} \text {does not, } \\
\kappa_{-} \in \operatorname{Card}^{\star} \text { if } K_{\lambda} \text { embeds in } \mathcal{G} \text { for all cardinal } \lambda<\kappa \text { and } K_{\kappa} \text { does not. }
\end{array}\right.
$$

Example. Let $\mathcal{G}$ be the disjoint union of graphs $\cup_{n<\aleph_{0}} K_{n}$ :


By definition, we have $C(\mathcal{G})=\aleph_{0-}$.
In addition to the union of graphs, we want to consider another natural operation:
Definition 1.28. We define the lexicographic product of graphs $\mathcal{G}$ and $\mathcal{F}$ as the graph $\mathcal{G}[\mathcal{F}]$ with set of vertices $G \times F$ and a symmetric relation given by:

$$
\left(g_{0}, f_{0}\right) R^{\mathcal{G}[\mathcal{F}]}\left(g_{1}, f_{1}\right) \Leftrightarrow \begin{cases}g_{0}=g_{1} & \text { and } f_{0} R^{\mathcal{F}} f_{1} \\ g_{0} \neq g_{1} & \text { and } g_{0} R^{\mathcal{G}} g_{1}\end{cases}
$$

Example. Consider the lexicographic product of $K_{4}$ and $K_{3}$. We simply obtain $K_{12}$ :


If $C(\mathcal{G}), C(\mathcal{F}) \in$ Card are cardinals greater or equal to 2 , we have by pigeonhole principle that

$$
C(\mathcal{G}[\mathcal{F}])=C(\mathcal{G}) \times C(\mathcal{F})
$$

This gives us the intuition of how one can define the product of cardinals in Card*. Let us look at two examples:
Example. Consider $\mathcal{G}=\cup_{n<\omega} K_{\aleph_{n}}$ and $\mathcal{F}=\cup_{\alpha<\omega_{1}} K_{\aleph_{\alpha}}$. Then we have $C(\mathcal{G})=\aleph_{\omega-}$ and $C(\mathcal{F})=\aleph_{\omega_{1}-}$. If we consider the lexicographic product with $K_{\aleph_{0}}$, we obtain:

- $C\left(K_{\aleph_{0}}[\mathcal{G}]\right)=\aleph_{\omega}$,
- $C\left(K_{\aleph_{0}}[\mathcal{F}]\right)=\aleph_{\omega_{1}-}$.

We leave the proof to the reader, with the following picture for the intuition:


As a consequence, one might be tempted to write $\aleph_{0} \cdot \aleph_{\omega-}=\aleph_{\omega}$ and $\aleph_{0} \cdot \aleph_{\omega_{1}-}=\aleph_{\omega_{1}-}$. This is what we want to define now.

Arithmetic on Card ${ }^{\star}$ We first define the cofinality of a cardinal $\lambda$ in Card $^{\star}$ as the cofinality of $\operatorname{act}(\lambda)$, denoted by $c f(\lambda)$. Secondly, we define the following operations:
Definition 1.29. Let $\Lambda=\left(\lambda_{i}\right)_{i \in I}$ be a sequence in $\operatorname{Card}^{\star}$. Let $\lambda=\sup _{i \in I}\left(\operatorname{act}\left(\lambda_{i}\right)\right) \in$ Card be the supremum in the usual sense, and $\operatorname{supp}(\Lambda)=\left\{i \in I \mid \lambda_{i} \neq 0\right\}$. We find a partition $I_{1} \cup I_{2} \cup I_{3}$ of $I$ such that:

$$
\Lambda=\left(\lambda_{i}\right)_{i \in I_{1}} \cup\left(\lambda_{-}\right)_{i \in I_{2}} \cup(\lambda)_{i \in I_{3}},
$$

where $\lambda_{i}<\lambda_{-}$for $i \in I_{1}$. We define sup ${ }^{\star}$ as follows:
$-\sup _{i \in I}^{\star}\left(\lambda_{i}\right)=\left\{\begin{array}{l}\lambda \text { if }\left|I_{3}\right| \neq \emptyset . \\ \lambda_{-} \text {otherwise. }\end{array}\right.$
If $|I|$ and $\lambda$ are finite, the definition of the sum $\sum^{\star}$ in Card $^{\star}$ is the sum in the usual sense:

$$
\sum_{i \in I}^{\star} \lambda_{i}=\sum_{i \in I} \operatorname{act}\left(\lambda_{i}\right)=\sum_{i \in I} \lambda_{i}
$$

Otherwise, we set:

- $\sum_{i \in I}^{\star} \lambda_{i}=\left\{\begin{array}{l}|\operatorname{supp}(\Lambda)| \text { if }|\operatorname{supp}(\Lambda)| \geq \lambda, \\ \left\{\begin{array}{l}\lambda \text { if } I_{3} \neq \emptyset, \\ \lambda \text { if }\left|I_{2}\right| \geq \operatorname{cf}(\lambda), \\ \lambda \text { if } \sup _{i \in I_{1}}\left(\operatorname{act}\left(\lambda_{i}\right)\right)=\lambda, \\ \lambda_{-} \text {otherwise. }\end{array} \quad \text { if }|\operatorname{supp}(\Lambda)|<\lambda .\right.\end{array}\right.$

For $\lambda, \mu \in \operatorname{Card}, \lambda$ limit cardinal, we define the product ${ }^{\star}:$ Card $\times \operatorname{Card}^{\star} \rightarrow \operatorname{Card}^{\star}$ in terms of sum:

- $\sum_{\mu}^{\star} \lambda_{-}=\mu \cdot{ }^{\star} \lambda_{-}=\left\{\begin{array}{l}\lambda_{-} \text {if } \mu<\operatorname{cf}(\lambda), \\ \mu \cdot \lambda \text { if } \mu \geq \operatorname{cf}(\lambda) .\end{array}\right.$

We see in particular that, under these definition, sup ${ }^{\star}$ and $\sum^{\star}$ do not necessary coincide anymore when there are infinite. However, it is clear that we recover the usual definition via the actualisation map, as we have the following commutative diagrams:


Here are the promised examples:
Examples. - Consider the sequence $\Lambda_{1}=\aleph_{\omega-}, 1,2,3, \ldots$ We have sup ${ }^{\star} \Lambda_{1}=\sum^{\star} \Lambda_{1}=\aleph_{\omega-}$.

- Consider the sequences $\Lambda_{2}=\left(\aleph_{\omega-}\right)_{i<\omega}=\aleph_{\omega-}, \aleph_{\omega-}, \ldots$ and $\Lambda_{3}=\left(\aleph_{i}\right)_{i<\omega}=\aleph_{0}, \aleph_{1}, \ldots$. We have $\sup ^{\star} \Lambda_{2}=\sup ^{\star} \Lambda_{3}=\aleph_{\omega-}$ and $\sum^{\star} \Lambda_{2}=\sum^{\star} \Lambda_{3}=\aleph_{\omega}$.
- Consider $\Lambda_{4}=\left(\aleph_{i}\right)_{i<\omega} \cup\left(\aleph_{2 \omega-}\right)$. Then $\sup ^{\star} \Lambda_{4}=\sum^{\star} \Lambda_{4}=\aleph_{2 \omega-}$.
- We have $\aleph_{0} \cdot \aleph_{\omega-}=\aleph_{\omega}, \aleph_{0} \cdot \aleph_{\omega_{1}-}=\aleph_{\omega_{1}-}$ and $\aleph_{1} \cdot \aleph_{\omega_{1}-}=\aleph_{\omega_{1}}$.

Now, we go back to the burden.
Burden, strength and Card ${ }^{\star}$ In Definition 1.13, the burden of the complete theory $T$ is the supremum (in Card $\cup\{\infty\}$ ) of depth of inp-patterns in $T$. However this supremum is not necessarily attained by an actual inp-pattern. This distinction is in particular motivated by the following definition:

Definition 1.30 ([1]). A complete theory is called strong if there is no inp-pattern of infinite depth in $T$.

One sees that, paradoxically, some strong theories have burden $\aleph_{0}$ and some theories of burden $\aleph_{0}$ are not strong (see examples below). In other words, the definition of burden we gave failed to characterize strength. We will use Adler's convention (see [1]) which gives a solution to this problem: the burden takes value in Card ${ }^{\star} \cup\{\infty\}$. We will indicate when we use this convention by writing bdn ${ }^{\star}$ instead of bdn.

Definition 1.31. (second definition of burden) Let $T$ be a complete theory. We denote by $\mathcal{S}$ the set of sorts.

- The burden $\operatorname{bdn}^{\star}(\pi(x))$ of a partial type $\pi(x)$ is the supremum in $\operatorname{Card}^{\star} \cup\{\infty\}$ of the depths of inp-patterns in $p(x)$.
- The cardinal $\sup _{S \in \mathcal{S}}^{\star} \operatorname{bdn}^{\star}\left(\left\{x_{S}=x_{S}\right\}\right)$ where $x_{S}$ is a single variable from the sort $S$, is called the burden of the theory $T$, and it is denoted by $\kappa_{\text {inp }}^{1 \star}(T)$ or by $\operatorname{bdn}^{\star}(T)$.

In other words, if the supremum $\lambda \in$ Card of depth of inp-patterns is attained, the burden is equal to $\lambda$. Otherwise, it is equal to $\lambda_{-}$. In particular, strong theories are exactly theories of burden at most $\aleph_{0-}$. One can check that every lemma in the previous paragraph -and its proof- still hold. Let us give a formal definition:

Definition 1.32. Let $\mathcal{M}_{i}=\left(M_{i}, \ldots\right)$ be a structure in a language $\mathrm{L}_{i}$, for $i \in I$ a set of indices. We define the following multisorted structure:

- The disjoint union

$$
\bigcup_{i} \mathcal{M}_{i}=\left\{\left(M_{i}, \ldots\right)\right\}_{i \in I}
$$

with a sort for each $\mathcal{M}_{i}$ 's.

- The direct product

$$
\prod_{i \in I} \mathcal{M}_{i}=\left\{\prod_{i \in I} M_{i},\left(M_{i}, \ldots\right)_{i \in I},\left(\pi_{i}: \prod_{j \in I} M_{j} \rightarrow M_{i}\right)_{i \in I}\right\}
$$

with a sort for each $\mathcal{M}_{i}$ 's and a sort for the product and where $\pi_{i}: \prod_{j \in I} M_{j} \rightarrow M_{i}$ is the natural projection.

We have the following fact:
Fact 1.33. - The sorts $\mathcal{M}_{i}$ in the union $\cup_{i \in I} \mathcal{M}_{i}$ are stably embedded and pairwise orthogonal.


- The direct product $\prod_{i \in I} \mathcal{M}_{i}$ eliminates quantifiers relative to the sorts $\mathcal{M}_{i}$. In particular, the sorts $\mathcal{M}_{i}$ are stably embedded, and pairwise orthogonal.


Proof. The first point is easy and is solved by simple inspection on formulas. For the second point, we leave to the reader to prove quantifier elimination.

Stable embeddedness is clear by inspection: a formula $\phi\left(x_{i}, a,\left(a_{j}\right)_{j \in J}\right)$ with variable $x_{i} \in \mathcal{M}_{i}$ and parameters $a \in \mathcal{M}$ and $a_{j} \in \mathcal{M}_{j}$ for $j \in I$ and without $M$-sorted quantifiers is a disjunction of formulas of the form

$$
\phi_{i}\left(x_{i}, a_{i}, \pi_{i}(a)\right) \wedge \phi(a) \wedge \bigwedge_{j \in I \backslash\{i\}} \phi_{j}\left(a_{j}, \pi_{j}(a)\right),
$$

where $\phi_{i}\left(x_{i}, a_{i}, \pi_{i}(a)\right)$ is an $\mathrm{L}_{i}$-formula, $\phi_{j}\left(a_{j}, \pi_{j}(a)\right)$ are closed $\mathrm{L}_{j}$-formula $j \in I \backslash\{i\}$ and $\phi(a)$ is a closed formula in the empty language. It is clearly equivalent to an $L_{i}$-formula with parameters in $\mathcal{M}_{i}$, as a closed formula is true or false and can be replaced either by $x_{i}=x_{i}$ or by $x_{i} \neq x_{i}$. Orthogonality is also clear: for the same reason, a formula $\phi\left(x_{i}, x_{j}, a,\left(a_{j}\right)_{j \in J}\right)$ with variable $x_{i} \in \mathcal{M}_{i}$ and $x_{j} \in \mathcal{M}_{j}$ without $M$-sorted quantifiers is equivalent to a disjunction of formulas of the form

$$
\phi_{i}\left(x_{i}, a_{i}, \pi_{i}(a)\right) \wedge \phi_{j}\left(x_{j}, a_{j}, \pi_{j}(a)\right)
$$

Naturally, we have a generalisation of Proposition 1.24 for infinite products.
Proposition 1.34. Let $\mathcal{M}_{i}=\left(M_{i}, \ldots\right)$ be a structure in a language $\mathrm{L}_{i}$, for $i \in I$ a set of indices. Assume they are not all finite. One has:

- $\operatorname{bdn}^{\star}\left(\bigcup_{i \in I} \mathcal{M}_{i}\right)=\sup _{i \in I}^{\star} \operatorname{bdn}^{\star}\left(M_{i}\right)$,
- $\operatorname{bdn}^{\star}\left(\prod_{i \in I} \mathcal{M}_{i}\right)=\sum_{i \in I}^{\star} \operatorname{bdn}^{\star}\left(M_{i}\right)$.

Remark 1.35. If all structures $\mathcal{M}_{i}$ are finite, there are two cases: either $\#\left\{i \in I\left|\left|M_{I}\right|>1\right\}\right.$ is infinite and $\operatorname{bdn}^{\star}\left(\prod_{i \in I} \mathcal{M}_{i}\right)=1$, or $\#\left\{i \in I\left|\left|M_{I}\right|>1\right\}\right.$ is finite and $\operatorname{bdn}^{\star}\left(\prod_{i \in I} \mathcal{M}_{i}\right)=0$. As the condition $\left|M_{i}\right|>1$ cannot be seen in terms of burden, we must treat this case separately.
Proof. The first point is clear: an inp-pattern $P(x)$ in $\bigcup_{i} \mathcal{M}_{i}$ has to 'choose' in which sort $\mathcal{M}_{i}$ its variable $x$ lives. This sort, say $\mathcal{M}_{i_{0}}$, is stably embedded by Fact 1.33. By Remark 1.22 , the depth of $P(x)$ is bounded by $\operatorname{bdn}\left(\mathcal{M}_{i_{0}}\right)$. Going to the supremum, one sees that definitions match:

$$
\operatorname{bdn}^{\star}\left(\bigcup_{i \in I} \mathcal{M}_{i}\right)=\sup _{i \in I}^{\star} \operatorname{bdn}^{\star}\left(M_{i}\right)
$$

The second point is more subtle: if $Q(x)$ is an inp-pattern of depth $\mu$ in $\prod_{i \in I} \mathcal{M}_{i}$, with the variable $x$ in the main sort, then the pattern refers to the sorts $\mathcal{M}_{i}$ simultaneously.

Claim 1. Assume that $\prod_{i \in I} \mathcal{M}_{i}$ admits an inp-pattern of depth $\mu$. Then there is an inp-pattern of depth $\mu$ in $\prod_{i \in I} \mathcal{M}_{i}$ of the following form:

$$
\left\{\phi_{\alpha}\left(\pi_{f(\alpha)}(x), y_{f(\alpha)}\right),\left(a_{f(\alpha), j}\right)_{j<\omega}\right\}_{\alpha<\mu},
$$

for some function $f: \mu \rightarrow I$ and where $\phi_{\alpha}\left(x_{f(\alpha)}, y_{f(\alpha)}\right)$ is a $\mathcal{M}_{f(\alpha)}$-formula. In other word, we may assume that a line $\alpha$ "mentions" only one structure $\mathcal{M}_{i}$.

Proof. Let us assume that $\prod \mathcal{M}_{i}$ admits an inp-pattern $Q(x)=\left\{\psi_{\alpha}\left(x, \bar{y}_{\alpha}\right),\left(\bar{a}_{\alpha, j}\right)_{j<\omega}, k_{\alpha}\right\}_{\alpha<\mu}$ of depth $\mu \geq 2$. We assume the array $\left(\bar{a}_{\alpha, j}\right)_{\alpha<\mu, j<\omega}$ to be mutually indiscernible. To simplify the notation, a generic line of $Q(x)$ is denoted by $\left\{\psi(x, \bar{y}),\left(\bar{a}_{j}\right)_{j<\omega}, k\right\}$ (we drop the index $\alpha$ ). By relative quantifier elimination and by Lemma 1.23, we may assume that formulas $\psi(x, \bar{y})$ in $Q(x)$ are of the form

$$
\bigwedge_{n<N} x \neq y_{n} \wedge x=y \wedge \bigwedge \phi_{i}\left(\pi_{i}(x), y_{i}\right),
$$

where $\phi_{i}\left(x, y_{i}\right)$ are $\mathrm{L}_{i}$-formulas, $N \in \mathbb{N}$ and where $\bar{y}=\left(y_{1}, \ldots, y_{N}, y\right) \cup\left(y_{i}\right)_{i \in I}$ and $\bar{a}_{j}=$ $\left(a_{1 j}, \ldots, a_{N j}, a_{j}\right) \cup\left(a_{i j}\right)_{i \in I}$ for $j<\omega$. If the atomic formula $x=y$ does occur, for example in the first row, then consistency of paths contradicts $k_{2}$-inconsistency of the second row. Thus, we knows that formulas in $Q(x)$ are of the form

$$
\bigwedge_{n<N} x \neq y_{n} \wedge \bigwedge \phi_{i}\left(\pi_{i}(x), y_{i}\right)
$$

Now, the formula $\bigwedge_{n<N} x \neq y_{n}$ is co-finite. This implies that

$$
\left\{\bigwedge \phi_{i}\left(\pi_{i}(x), a_{i, j}\right)\right\}_{j<\omega}
$$

is $k+1$-inconsistent. Indeed, otherwise, for one (equivalently for all) $k+1$-increasing tuple $j_{0}<\cdots<$ $j_{k}<\omega$, the set

$$
\left\{\bigwedge \phi_{i}\left(\pi_{i}(x), a_{i, j_{1}}\right), \ldots, \bigwedge \phi_{i}\left(\pi_{i}(x), a_{i, j_{k}}\right)\right\}
$$

is satisfied by $a_{n, j_{l}}$ for some $n<N$ and $l \leq k$. Without loss of generality, assume that $n=N-1$ and $l=k$. Then, by mutual indiscernibility, $\left(a_{N-1, j}\right)_{j \geq k}$ are solutions of

$$
\left\{\bigwedge \phi_{i}\left(\pi_{i}(x), a_{i, 0}\right), \ldots, \bigwedge \phi_{i}\left(\pi_{i}(x), a_{i, k-1}\right)\right\}
$$

This contradicts the $k$-inconsistency of the line

$$
\left\{\bigwedge_{n<N} x \neq a_{n, j} \wedge \bigwedge \phi_{i}\left(\pi_{i}(x), a_{i, j}\right)\right\}_{j<\omega}
$$

unless $\left(a_{N-1, j}\right)_{j<\omega}$ is constant. In that case, this parameter can be ignore: replace the formula by

$$
\bigwedge_{n<N-1} x \neq y_{n} \wedge \bigwedge \phi_{i}\left(\pi_{i}(x), y_{i}\right)
$$

and we still have an inp-pattern. We get our contradiction by induction on $N$. Hence, we may assume that formulas $\psi(x, \bar{y})$ in $Q(x)$ are of the form

$$
\bigwedge \phi_{i}\left(\pi_{i}(x), y_{i}\right)
$$

We may now conclude using mutual indiscernibility that for at least one $i=: f(\alpha)$, the set

$$
\left\{\phi_{i}\left(\pi_{i}(x), a_{i, j}\right)\right\}_{j<\omega}
$$

is $k$-inconsistant. We may replace the formula $\bigwedge \phi_{i}\left(\pi_{i}(x), y_{i}\right)$ by $\phi_{f(\alpha)}\left(\pi_{f(\alpha)}, y_{f(\alpha)}\right)$. In other word, we may assume that only the index $i=f(\alpha)$ occurs in the formula of the line $\alpha$. We found an inp-pattern of the desired form.

We denote $\operatorname{bdn}^{\star}\left(\mathcal{M}_{i}\right)$ by $\lambda_{i}$ and $\sup _{i \in I}$ act $\left(\lambda_{i}\right) \in$ Card by $\lambda$. One immediate corollary is that

$$
\operatorname{bdn}^{\star}\left(\prod_{i \in I} \mathcal{M}_{i}\right)=\operatorname{bdn}^{\star}\left(\prod_{\substack{i \in I \\ \lambda_{i} \neq 0}} \mathcal{M}_{i}\right) \geq 1
$$

(notice that we used that some $\mathcal{M}_{i}$ is infinite). We may assume that $I=\operatorname{supp}\left(\lambda_{i}\right)_{i \in I}$. Now, the proof is straight forward and is just a case study. We distinguish six cases:

First case: the cardinals $|I|$ and $\lambda$ are finite. Then, this is immediate from the previous claim: $\operatorname{bdn}^{\star}\left(\prod \mathcal{M}_{i}\right)=\sum \lambda_{i}=\sum^{\star} \lambda_{i}$.

Second case: we have $|I| \geq \lambda$ and $|I| \geq \aleph_{0}$. Then, let $\left(b_{i, j}\right)_{j<\omega}$ be a sequence of pairwise distinct elements of $\mathcal{M}_{i}$. Let $x$ be a variable in the main sort. Then, $\left\{\pi_{i}(x)=y_{i},\left(b_{i, j}\right)_{j<\omega}\right\}_{i \in I}$ is an inp-pattern of depth $|I|$. We have $\operatorname{bdn}^{\star}\left(\prod M_{i}\right) \geq|I|$. Reciprocally, assume $\prod \mathcal{M}_{i}$ admits an inp-pattern $Q(x)$ of depth $\mu>|I|$. By the previous claim and pigeonhole principle, we find an inp-pattern of depth $\mu$ in some $\mathcal{M}_{i}$, which is a contradiction with $\lambda \leq|I|<\mu$. We get $\operatorname{bdn}^{\star}\left(\prod \mathcal{M}_{i}\right)=\sum_{i \in I}^{\star} \lambda_{i}=|I|$.

Third case: we have $|I|<\lambda$ and $\lambda_{i}=\lambda \geq \aleph_{0}$ for some $i \in I$. Then clearly $\operatorname{bdn}^{\star}\left(\Pi \mathcal{M}_{i}\right) \geq \lambda$. Again, by pigeonhole principle, one gets $\operatorname{bdn}^{\star}\left(\Pi \mathcal{M}_{i}\right) \leq \lambda$.

Fourth case: we have $|I|<\lambda$ and $\operatorname{cf}(\lambda) \leq \#\left\{i \in I \mid \lambda_{i}=\lambda_{-}\right\}$. Then, choose any sequence of cardinals $\left(\mu_{\alpha}\right)_{\alpha<\operatorname{cf}(\lambda)}$ with supremum $\lambda$ (in the usual sense) and $\mu_{\alpha}<\lambda$ for all $\alpha$. We can assume that $I=\operatorname{cf}(\lambda)$ and that we have an inp-pattern $Q_{i}\left(x_{i}\right)$ in $\mathcal{M}_{i}$ of depth $\mu_{i}$. The inp-pattern $Q(x)=$ $\cup_{i \in I} Q_{i}\left(\pi_{i}(x)\right)$ is of depth $\lambda$. We get $\operatorname{bdn}^{\star}\left(\prod \mathcal{M}_{i}\right)=\sum_{i \in I}{ }^{\star} \lambda_{i}=\lambda$.

Fifth case: we have $\sup \left\{\operatorname{act}\left(\lambda_{i}\right) \mid \lambda_{i} \notin\left\{\lambda_{-}, \lambda\right\}\right\}=\lambda$. We conclude as in the previous case that $\operatorname{bdn}^{\star}\left(\prod \mathcal{M}_{i}\right)=\sum_{i \in I} \lambda_{i}=\lambda$.

Last case: we are not in the above cases. Then, by the previous claim, there is no inp-pattern of depth $\lambda$ in $\prod \mathcal{M}_{i}$. We have then

$$
\operatorname{bdn}^{\star}\left(\prod \mathcal{M}_{i}\right)=\sum_{i \in I}^{\star} \lambda_{i}=\lambda_{-} .
$$

In a supersimple theory, the burden of a complete type is always finite (see [1]). Hence, supersimple theories are examples of strong theories.
Example. The following structures have burden $\aleph_{0-}$ :

- Any union structure $M=\bigcup_{n} M_{n}$, where for every $n \in \mathbb{N}, M_{n}$ is a structure of burden $n$.
- Any model of ACFA, the model companion of the theory of algebraically closed fields with an automorphism.
- Any model of $\mathrm{DCF}_{0}$, the theory of differentially closed fields.

The first example is clear by the previous discussion but could look artificial. It will naturally appear when we will discuss the burden of the $\mathrm{RV}_{<\omega}$-sort in mixed characteristic (see Section 3.3). The fact that the last two examples are of burden $\aleph_{0-}$ (and not finite) follows from the fact they are super-simple and from the next remark. Indeed, Chernikov and Hils noticed that, since such fields are infinite dimensional vector spaces over respectively their fixed field and their constant field, they must admit a pattern of size $n$ for every integer $n$ :

Remark 1.36. [9, Remark 5.3] Let $T$ be a simple theory and assume there is a $n$-dimensional typedefinable vector space $V$ over a type-definable infinite field $F$. Then there is a type in $V$ of burden $\geq n$.

Let us look at one natural example of a non-strong theory:

Remark 1.37. Let $k$ be an imperfect field of characteristic $p$, considered as a structure in the language of fields. Then $\operatorname{bdn}^{\star}(k) \geq \aleph_{0}$.

Proof. Let $e_{0}, e_{1} \in k$ be two linearly independent elements over $k^{p}$. Then, we have the definable injective map

$$
\begin{array}{rllc}
f_{2}: & k \times k & \rightarrow & k \\
& (a, b) & \mapsto & a^{p} e_{0}+b^{p} e_{1}
\end{array}
$$

By induction, we define for $n \geq 2$ :

$$
\begin{array}{cccc}
f_{n+1}: & k^{(n+1)} & \rightarrow & k \\
\left(a_{0}, \ldots, a_{n-1}\right) & \mapsto & f_{n}\left(a_{0}, \ldots, a_{n-3}, f_{2}\left(a_{n-2}, a_{n-1}\right)\right) .
\end{array}
$$

Then, consider the formula for $n \geq 0$ :

$$
\phi_{n}\left(x, y_{n}\right) \equiv \exists y_{0}, \ldots, y_{n-1}, y_{n+1} x=f_{n+2}\left(y_{0}, \ldots, y_{n+1}\right),
$$

and pairwise distinct parameters $b_{n, j} \in k$, for $j<\omega$. Then

$$
\left\{\phi_{n}\left(x, y_{n}\right),\left(b_{n, j}\right)_{j<\omega}\right\}_{n<\omega}
$$

is an inp-pattern of depth $\aleph_{0}$ : inconsistency of the rows is clear, and it is easy to show consistency of paths by compactness.

Remark 1.38. Similarly, one can show that if a model $\mathcal{M}$ is bi-interpretable on a unary set with the direct product $\mathcal{M} \times \mathcal{M}=\left(M \times M, \mathcal{M}, \pi_{1}: M \times M \rightarrow M, \pi_{2}: M \times M \rightarrow M\right)$, then $\operatorname{bdn}^{\star}(\mathcal{M}) \geq \aleph_{0}$ and it is never of the form $\lambda_{-}$where $\lambda$ is a cardinal of cofinality $\aleph_{0}$. Indeed, let us sketch an argument: Assume $\operatorname{bdn}^{\star}(\mathcal{M}) \geq \lambda_{-}$, let $\left(\lambda_{n}\right)_{n}$ be a sequence of cardinals cofinal in $\lambda$. We must find an inp-pattern of depth $\lambda$ in $\mathcal{M}$. Consider $f_{2}: D_{2} \longrightarrow \mathcal{M} \times \mathcal{M}$ an interpretation map where $D_{2}$ is a definable unary subset of $\mathcal{M}$ (and the relation $f(a)=f\left(a^{\prime}\right)$ is a definabe equivalence relation on $D_{2}$ ). We may define by induction an interpretation of $\prod_{n} \mathcal{M}$ on a unary set $D_{n} \subset \mathcal{M}$ by 'duplicating' the last component at each step:

$$
\begin{aligned}
& D_{2} \xrightarrow{f_{2}} \mathcal{M} \times \mathcal{M}, \\
& D_{3} \xrightarrow{f_{2}} \mathcal{M} \times D_{2} \xrightarrow{i d \times f_{2}} \mathcal{M} \times(\mathcal{M} \times \mathcal{M}), \\
& D_{4} \xrightarrow{f_{2}} \mathcal{M} \times D_{3} \xrightarrow{i d \times f_{2}} \mathcal{M} \times\left(\mathcal{M} \times D_{2}\right) \xrightarrow{i d \times i d \times f_{2}} \mathcal{M} \times(\mathcal{M} \times(\mathcal{M} \times \mathcal{M})),
\end{aligned}
$$

etc.

Now, given an inp-pattern $P_{n}(x)=\left\{\phi_{i}\left(x_{n}, y_{i}\right),\left(a_{i, j}\right)_{j<\omega}\right\}_{i<\lambda_{n}}$ in $\mathcal{M}$ of depth $\lambda_{n}$, we define a new pattern $P_{n}^{\prime}(x):=\left\{\phi_{i}^{\prime}\left(x, y_{i}^{\prime}\right),\left(a_{i, j}^{\prime}\right)_{j<\omega}\right\}_{i<\lambda_{n}}$ by taking the (definable) pre-image $\phi_{i}^{\prime}\left(D_{n+1}, a_{i, j}^{\prime}\right)$ of the definable set $\phi_{i}\left(M, a_{i, j}\right)$ via the map $\pi_{n} \circ f_{n+1}: D_{n+1} \rightarrow \mathcal{M}$ (where $\pi_{n}$ is the projection to the $n^{\text {th }}$ coordinate). We see then that $P_{n}^{\prime}(x)$ is also an inp-pattern. Furthermore, define $P^{\prime}(x)=\bigcup_{n} P_{n}^{\prime}(x)$. One can see, using compactness, that all paths are consistent, and that $P^{\prime}(x)$ is an inp-pattern of depth $\lambda$ in $\mathcal{M}$.

### 1.2 On model theory of algebraic structures

### 1.2.1 Valued fields

We gather here some facts on valued fields. After some general statement on indiscernible sequences, we will introduce the RV-sort. Then, we will list the theories of valued fields that we will consider. We will recall some basics of the theory of Kaplansky. The reader can refer to [24] for more details.

We will need a few lemmas, such as a kind of transitivity of pseudo-limits, and a case study of indiscernible sequences. A valued field will be typically denoted by $\mathcal{K}=(K, \Gamma, k$, val) where $K$ is the field (main sort), $\Gamma$ the value group and $k$ the residue field. The valuation is denoted by val, the maximal ideal $\mathfrak{m}$ and the valuation ring $\mathcal{O}$. We recall the two traditional languages of valued fields.

## Notation and languages

We will work in different (many-sorted) languages. Let us define two of them:

- $\mathrm{L}_{\text {div }}=\{K, 0,1,+, \cdot, \mid\}$, where $\mid$ is a binary relation symbol, interpreted by the division:

$$
\text { for } a, b \in K, a \mid b \text { if and only if } \operatorname{val}(a) \leq \operatorname{val}(b) .
$$

- $\mathrm{L}_{\Gamma, k}=\{K, 0,1,+, \cdot\} \cup\{k, 0,1,+, \cdot\} \cup\{\Gamma, 0, \infty,+,<\} \cup\left\{\mathrm{val}: K \rightarrow \Gamma\right.$, Res $\left.: K^{2} \rightarrow k\right\}$.
where Res : $K^{2} \rightarrow k$ is the two-place residue map, interpreted as follows:

$$
\operatorname{Res}(a, b)=\left\{\begin{array}{l}
\operatorname{res}(a / b) \text { if } \operatorname{val}(a) \geq \operatorname{val}(b) \neq \infty, \\
0 \text { otherwise } .
\end{array}\right.
$$

In the next paragraphs, we will also introduce the many-sorted languages $\mathrm{L}_{\mathrm{RV}}$ and $\mathrm{L}_{\mathrm{RV}}{ }_{<\omega}$ which involve the leading term structures $R V$ and $R V_{<\omega}$.

By bi-interpretability, a theory of valued fields can be expressed indifferently in either of these languages. Let $\mathcal{K}$ be a valued field. If the context is clear, we will often abusively denote by $K, \Gamma, k, \mathrm{RV}, \ldots$ the sorts in $\mathcal{K}$. In general, the sorts of a valued field $\mathcal{L}$ will be denoted by $L, \Gamma_{L}, k_{L}, \mathrm{RV}_{L} \ldots$ and of a valued field $\mathcal{K}^{\prime}$ by $K^{\prime}, \Gamma^{\prime}, k^{\prime}, \mathrm{RV}^{\prime}, \ldots$ etc.

Pseudo-Cauchy sequences We will discuss here some simple facts about mutually indiscernible arrays in a valued field $\mathcal{K}$. We will denote by $\overline{\mathbb{Z}}$ the set of integers with extreme elements $\{-\infty, \infty\}$. We recall the definition of pseudo-Cauchy sequences:
Definition 1.39. Let $(I,<)$ be a totally ordered index set without greatest element. A sequence $\left(a_{i}\right)_{i \in I}$ of elements of $K$ is pseudo-Cauchy if there is $i \in I$ such that for all indices $i<i_{1}<i_{2}<i_{3}$, $\operatorname{val}\left(a_{i_{2}}-a_{i_{1}}\right)<\operatorname{val}\left(a_{i_{3}}-a_{i_{2}}\right)$. We say that $a \in K$ is a pseudo limit of the pseudo-Cauchy sequence $\left(a_{i}\right)_{i \in I}$ and we write $\left(a_{i}\right)_{i \in I} \Rightarrow a$ if there is $i \in I$ such that for all indices $i<i_{1}<i_{2}$, we have $\operatorname{val}\left(a-a_{i_{1}}\right)=\operatorname{val}\left(a_{i_{2}}-a_{i_{1}}\right)$.

The next two lemmas give some useful properties of indiscernible pseudo-Cauchy sequences.
Lemma 1.40. 1. Assume $\left(a_{i}\right)_{i<\omega}$ is an indiscernible sequence and $a$ is a pseudo limit of $\left(a_{i}\right)_{i<\omega}$. Then for any $i, \operatorname{val}\left(a_{i}-a\right)=\operatorname{val}\left(a_{i}-a_{i+1}\right)$ depends only on $i$ and not on the chosen limit a (for general pseudo-Cauchy sequence, this holds only for $i$ big enough).
2. For three mutually indiscernible sequences $\left(a_{i}\right)_{i<\omega},\left(b_{i}\right)_{i<\omega}$ and $\left(c_{i}\right)_{i<\omega}$, if $\left(a_{i}\right)_{i<\omega} \Rightarrow b_{0}$ and $\left(b_{i}\right)_{i<\omega} \Rightarrow c_{0}$, then we have $\left(a_{i}\right)_{i<\omega} \Rightarrow c_{0}$.

3. If $\left(a_{i}\right)_{i \in \overline{\mathbb{Z}}}$ is an indiscernible sequence in $K$, then $\left(a_{i}\right)_{i \in \omega} \Rightarrow a_{\infty}$ or $\left(a_{-i}\right)_{i \in \omega} \Rightarrow a_{-\infty}$ or for $i \neq j$, $\operatorname{val}\left(a_{i}-a_{j}\right)$ is constant (in this last case, $\left(a_{i}\right)_{i \in \overline{\mathbb{Z}}}$ will be called a fan).
Proof. 1. By definition of a pseudo-Cauchy sequence, $\left(\operatorname{val}\left(a_{i}-a_{i+1}\right)\right)_{i}$ is eventually strictly increasing. By indiscernibility, it is strictly increasing. Let $i_{0}$ be such that $\operatorname{val}\left(a-a_{i}\right)=\operatorname{val}\left(a_{i+1}-a_{i}\right)$ for all $i>i_{0}$. Then $\operatorname{val}\left(a-a_{i_{0}}\right)=\min \left(\operatorname{val}\left(a-a_{i_{0}+1}\right), \operatorname{val}\left(a_{i_{0}+1}-a_{i_{0}}\right)\right)=\min \left(\operatorname{val}\left(a_{i_{0}+2}-\right.\right.$ $\left.\left.a_{i_{0}+1}\right), \operatorname{val}\left(a_{i_{0}+1}-a_{i_{0}}\right)\right)=\operatorname{val}\left(a_{i_{0}+1}-a_{i_{0}}\right)$. It holds also for $i=i_{0}$ and we can reiterate.
2. Notice that, by mutual indiscernibility and (1), $\operatorname{val}\left(a_{i}-b_{0}\right)=\operatorname{val}\left(a_{i}-a_{i+1}\right)=\operatorname{val}\left(a_{i}-b_{j}\right)$ for any $i, j<\omega$, i.e. $\left(a_{i}\right)_{i<\omega} \Rightarrow b_{j}$ for any $j$. Similarly, $\left(b_{i}\right)_{i<\omega} \Rightarrow c_{j}$ for any $j$. We have $\operatorname{val}\left(b_{0}-b_{1}\right) \geq$ $\operatorname{val}\left(b_{0}-a_{i}\right)=\operatorname{val}\left(a_{i}-b_{1}\right)$. If $\operatorname{val}\left(b_{0}-b_{1}\right)=\operatorname{val}\left(b_{0}-a_{i}\right)$, we have by mutual indiscernibility that $\left(\operatorname{val}\left(b_{0}-a_{i}\right)\right)_{i<\omega}$ is constant, which is a contradiction with $\left(a_{i}\right)_{i<\omega} \Rightarrow b_{0}$. Then, we have $\operatorname{val}\left(b_{0}-c_{0}\right)=\operatorname{val}\left(b_{0}-b_{1}\right)>\operatorname{val}\left(b_{0}-a_{i}\right)$. As $\operatorname{val}\left(a_{i}-c_{0}\right) \geq \min \left(\operatorname{val}\left(a_{i}-b_{0}\right), \operatorname{val}\left(b_{0}-c_{0}\right)\right)$, we deduce that $\operatorname{val}\left(a_{i}-c_{0}\right)=\operatorname{val}\left(a_{i}-b_{0}\right)$ for all i, i.e. $\left(a_{i}\right) \Rightarrow c_{0}$.
3. It is immediate by indiscernibility (consider for example $\operatorname{val}\left(a_{0}-a_{1}\right)$ and $\operatorname{val}\left(a_{1}-a_{2}\right)$ ).

Lemma 1.41. Let $\left(a_{j}\right)_{j \in \mathbb{Z}}$ and $\left(b_{l}\right)_{l \in \mathbb{Z}}$ two mutually indiscernible sequences in $K$ such that $\left(\operatorname{val}\left(a_{j}-b_{l}\right)\right)_{j, l}$ is not constant. At least one of the following occurs:

1. $\left(a_{j}\right)_{j<\omega} \Rightarrow b_{0}$,
2. $\left(b_{l}\right)_{l<\omega} \Rightarrow a_{0}$,
3. $\left(a_{-j}\right)_{j<\omega} \Rightarrow b_{0}$,
4. $\left(b_{-l}\right)_{l<\omega} \Rightarrow a_{0}$.

Note that if for example $\left(b_{l}\right)_{l<\omega} \Rightarrow a_{0}$, then by mutual indiscernibility, $\left(b_{l}\right)_{l<\omega} \Rightarrow a_{j}$ for every $j \in \mathbb{Z}$.
Proof. Since val $\left(a_{j}-b_{l}\right)$ is not constant, using the mutual indiscernibility, one of the following occurs:

1. $\operatorname{val}\left(a_{0}-b_{0}\right)<\operatorname{val}\left(a_{1}-b_{0}\right)$,
2. $\operatorname{val}\left(a_{0}-b_{0}\right)<\operatorname{val}\left(a_{0}-b_{1}\right)$,
3. $\operatorname{val}\left(a_{0}-b_{0}\right)<\operatorname{val}\left(a_{-1}-b_{0}\right)$,
4. $\operatorname{val}\left(a_{0}-b_{0}\right)<\operatorname{val}\left(a_{0}-b_{-1}\right)$.

Indeed, if 1 . and 3. do not hold, then the sequence $\left(\operatorname{val}\left(b_{0}-a_{j}\right)\right)_{j \in \mathbb{Z}}$ is constant. If 2. and 4. do not hold, then the sequence $\left(\operatorname{val}\left(b_{l}-a_{0}\right)\right)_{l \in \mathbb{Z}}$ is constant. This cannot be true for both sequences as it would contradict the assumption. We conclude by indiscernibility.


The leading term structure We will now define the RV-sort (or RV-sorts, as we may need to consider more than one sort) -an intermediate structure between the valued field and its value group and residue field. We also introduce corresponding languages $\mathrm{L}_{\mathrm{RV}}$ and $\mathrm{L}_{\mathrm{RV}}<\omega$. This paragraph is largely inspired by [17], which one can use as a reference. Let $K$ be a Henselian valued field of characteristic $(0, p)$ with $p \geq 0$, of value group $\Gamma$ and residue field $k$. If $\delta \in \Gamma_{\geq 0}$, we denote by $\mathfrak{m}_{\delta}$ the ideal of the valuation ring $\mathcal{O}$ defined by $\{x \in \mathcal{O} \mid v(x)>\delta\}$. The leading term structure of order $\delta$ is the quotient group

$$
\mathrm{RV}_{\delta}^{\star}:=K^{\star} /\left(1+\mathfrak{m}_{\delta}\right) .
$$

The quotient map is denoted by $\mathrm{rv}_{\delta}: K^{\star} \rightarrow \mathrm{RV}_{\delta}^{\star}$. The valuation val $: K^{\star} \rightarrow \Gamma$ induces a group homomorphism $\operatorname{val}_{\mathrm{rv}_{\delta}}: \mathrm{RV}_{\delta}^{\star} \rightarrow \Gamma$. Since $\mathfrak{m}=\mathfrak{m}_{0}$ and $k^{\star}:=\left(\mathcal{O} / \mathfrak{m}_{0}\right)^{\star} \simeq \mathcal{O}^{\times} /\left(1+\mathfrak{m}_{0}\right)$, we have the following short exact sequence:

$$
1 \rightarrow k^{\star} \xrightarrow{\iota} \mathrm{RV}_{0}^{\star} \xrightarrow{\mathrm{val}_{\mathrm{rv}}} \mathrm{C} \rightarrow 0
$$

In general, we denote by $\mathcal{O}_{\delta}$ the ring $\mathcal{O} / m_{\delta}$, called the residue ring of order $\delta$. One has $\mathcal{O}_{\delta}^{\times} \simeq$ $\mathcal{O}^{\times} /\left(1+\mathfrak{m}_{\delta}\right)$ and the following exact sequence:

$$
1 \rightarrow \mathcal{O}_{\delta}^{\times} \xrightarrow{\iota_{\delta}} \mathrm{RV}_{\delta}^{\star} \xrightarrow{\mathrm{val}_{\mathrm{rv}}} \Gamma \rightarrow 0
$$

Furthermore, as $\mathfrak{m}_{\gamma} \subseteq \mathfrak{m}_{\delta}$ for any $\delta \leq \gamma$ in $\Gamma_{\geq 0}$, we have a projection map $\mathrm{RV}_{\gamma}^{\star} \rightarrow \mathrm{RV}_{\delta}^{\star}$ denoted by $\operatorname{rv}_{\gamma \rightarrow \delta}$ or simply by $\mathrm{rv}_{\delta}$. We add a new constant $\mathbf{0}$ to the sort $\mathrm{RV}_{\delta}^{\star}$ and we write $\mathrm{RV}_{\delta}:=\mathrm{RV}_{\delta}^{\star} \cup\{\mathbf{0}\}$. We set the following properties:

- for all $\mathbf{x} \in \mathrm{RV}_{\delta}, \mathbf{0} \cdot \mathbf{x}=\mathbf{x} \cdot \mathbf{0}=\mathbf{0}$.
- $\operatorname{val}_{\mathrm{rv}_{\delta}}(\mathbf{0})=\infty, \quad \operatorname{rv}_{\delta}(0)=\mathbf{0}$.

Proposition 1.42. For any $a, b \in K$ and $\delta \in \Gamma_{\geq 0}, \operatorname{rv}_{\delta}(a)=\operatorname{rv}_{\delta}(b)$ if and only if $\operatorname{val}(a-b)>\operatorname{val}(b)+\delta$ or $a=b=0$.

Proof. This follows easily from the definition: assume $\operatorname{rv}_{\delta}(a)=\operatorname{rv}_{\delta}(b)$ and $a \neq 0$. Then $a=b(1+\mu)$ for some $\mu \in m_{\delta}$ and $\operatorname{val}(a-b)=\operatorname{val}(b)+\operatorname{val}(\mu)>\operatorname{val}(b)+\delta$. Conversely, if $\operatorname{val}(a-b)>\operatorname{val}(b)+\delta$, one can write $a=b\left(1+\frac{(a-b)}{b}\right)$.

As a group quotient, the sort $\mathrm{RV}_{\delta}$ is endowed with a multiplication. As we will see, it also inherits from the field some kind of addition.
Notation. Let $0 \leq \delta_{1}, \delta_{2}, \delta_{3}$ be three elements of $\Gamma$ and $\mathbf{x} \in \mathrm{RV}_{\delta_{1}}, \mathbf{y} \in \mathrm{RV}_{\delta_{2}}, \mathbf{z} \in \mathrm{RV}_{\delta_{3}}$ three variables. Then we define the following formulas:

$$
\oplus_{\delta_{1}, \delta_{2}, \delta_{3}}(\mathbf{x}, \mathbf{y}, \mathbf{z}) \equiv \quad \exists a, b \in K \operatorname{rv}_{\delta_{1}}(a)=\mathbf{x} \wedge \operatorname{rv}_{\delta_{2}}(b)=\mathbf{y} \wedge \operatorname{rv}_{\delta_{3}}(a+b)=\mathbf{z}
$$

In our study of valued fields, we will consider the structures $R V$ and $R V<\omega$, that we define now:
Definition 1.43. The RV-sort of a valued field $\mathcal{K}$ is the first order leading term structure

$$
\mathrm{RV}=\left(\mathrm{RV}_{0}, \cdot, \oplus_{0,0,0}, \mathbf{0}, \mathbf{1}\right)
$$

endowed with its natural structure of abelian group and the ternary predicate described above. Following the usual convention, we drop the index 0 , and write $\mathrm{RV}, \oplus$ and rv instead of $\mathrm{RV}_{0}, \oplus_{0,0,0}$ and $\mathrm{rv}_{0}$.

Fact 1.44 ( Flenner, [17, Proposition 2.8]). The three-sorted structure $\left\{(\mathrm{RV}, \mathbf{1}, \cdot, \mathbf{0}, \mathbf{1}),(k, 0,1,+, \cdot),(\Gamma, 0,+,<), \iota, \mathrm{val}_{\mathrm{rv}}\right\}$ and the one-sorted structure $\{\mathrm{RV}, \mathbf{0}, \cdot, \oplus\}$ are bi-interpretable on unary sets.

This will mean, in the context of this paper, that these two points of view are equivalent, and we will swap between one to the other indifferently (see Fact 1.18).

We defined the leading term language $\mathrm{L}_{\mathrm{RV}}$ as the multisorted language with

- a sort for $K$ and RV.
- the ring language for $K$,
- the (multiplicative) group language as well as the symbol 0 for RV.
- the ternary relation symbol $\oplus$ and the function symbols rv.

The structure $\mathcal{K}=(K, R V, r v)$ becomes a structure in this language where all symbols are interpreted as before. This language is also bi-interpretable (without parameters) with the usual languages of valued fields, e.g. with $\mathrm{L}_{\text {div }}$ (see [17, Proposition 2.8]).

Also, notice that the symbol $\oplus$ suggests a binary operation. Occasionally, we will indeed write $\operatorname{rv}(a) \oplus \operatorname{rv}(b)$ for $a, b \in K$ to denote the following element:

$$
\operatorname{rv}(a) \oplus \operatorname{rv}(b):= \begin{cases}\operatorname{rv}(a+b) & \text { if } \operatorname{val}(a+b)=\min (\operatorname{val}(a), \operatorname{val}(b)) \\ \mathbf{0} & \text { otherwise }\end{cases}
$$

It is not hard to see that this is independent of the choice of representatives of $\operatorname{rv}(a)$ and $\operatorname{rv}(b)$. We will write $\bigoplus_{i \in I} \mathbf{a}_{i}$ for $I$ a set of indices and $\mathbf{a}_{i} \in \mathrm{RV}$, when such a sum does not depend on any choices of parentheses. Notice that even if we denote this law $\oplus$ with an addition symbol, this operation is not associative. Also, if $a, b \in K$ with $\operatorname{val}(a)<\operatorname{val}(b)$, we have that $\operatorname{rv}(a) \oplus \operatorname{rv}(b)=\operatorname{rv}(a)=\operatorname{rv}(a+b)$. It is not true in general that $\operatorname{rv}(a+b)=\operatorname{rv}(a) \oplus \operatorname{rv}(b)$ (choose $a, b \in K$ such that $\operatorname{rv}(a)=-\operatorname{rv}(b)$ and $a \neq-b)$. When we have indeed that $\mathrm{rv}(a+b)=\operatorname{rv}(a) \oplus \operatorname{rv}(b)$, we say that the sum $\operatorname{rv}(a) \oplus \operatorname{rv}(b)$ is well-defined.

In the specific context of mixed characteristic Henselian valued fields, we might have to consider a larger structure:

Definition 1.45. Assume that $\mathcal{K}$ is a valued field of characteristic 0 and residue characteristic $p \geq 0$. We reserve now the notation $\delta_{n}$ for $\delta_{n}=\operatorname{val}\left(p^{n}\right)$. We write $\mathrm{RV}<\omega$ for the union of sorts leading term structure of finite order $\left\{\left(\mathrm{RV}_{\delta_{n}}\right)_{n<\omega},\left(\oplus_{\delta_{l}, \delta_{m}, \delta_{n}}\right)_{n<l, m},\left(\operatorname{rv}_{\delta_{n} \rightarrow \delta_{m}}\right)_{m<n<\omega}\right\}$ endowed with ternary predicates $\oplus_{\delta_{l}, \delta_{m}, \delta_{n}}$ and a projective system of maps $\left(\operatorname{rv}_{\delta_{n} \rightarrow \delta_{m}}\right)_{m<n<\omega}$. We also write val $\mathrm{rv}_{<\omega}: \mathrm{RV}_{<\omega} \rightarrow$ $\Gamma \cup\{\infty\}$ for $\bigcup_{n<\omega}\left(\operatorname{val}_{\mathrm{rv}_{\delta_{n}}}: \operatorname{RV}_{\delta_{n}} \rightarrow \Gamma \cup\{\infty\}\right)$, etc.

Remark 1.46. In equicharacteristic 0 , we have that $\delta_{n}:=\operatorname{val}\left(p^{n}\right)=0$ for all $n<\omega$. This leads to identifying $\bigcup_{n<\omega} \mathrm{RV}_{\delta_{n}}$ with $\mathrm{RV}=\mathrm{RV}_{0}$.

In Section 3.3, we will have to use another language to describe the induced structure on $\mathrm{RV}_{<\omega}$ :
Fact 1.47 ( Flenner, [17, Proposition 2.8]). The structure

$$
\left\{\left(\mathrm{RV}_{\delta_{n}}\right)_{n<\omega},\left(\oplus_{\delta_{l}, \delta_{m}, \delta_{n}}\right)_{n<l, m},\left(\mathrm{rv}_{\delta_{n} \rightarrow \delta_{m}}\right)_{m<n<\omega}\right\}
$$

and the structure

$$
\left.\left.\begin{array}{r}
\left\{\left(\operatorname{RV}_{\delta_{n}}\right)_{n<\omega},\left(\mathcal{O}_{\delta_{n}}, \cdot,+, 0,1\right)_{n<\omega},(\Gamma,+, 0,<),\left(\operatorname{val}_{\mathrm{rv}}^{\delta_{n}}\right.\right.
\end{array}\right)_{n<\omega}, ~ 子, ~\left(\mathcal{O}_{\delta_{n}}^{\times} \rightarrow \operatorname{RV}_{\delta_{n}}^{\times}\right)_{n<\omega},\left(\operatorname{rv}_{\delta_{n} \rightarrow \delta_{m}}\right)_{m<n<\omega}\right\}
$$

are bi-interpretable on unary sets.

As before, this is only to say that one can recover the valuation using the symbols $\oplus$ (see [17, Proposition 2.8]). And again, this fact means that we will be able to swap between one language to the other indifferently (see Fact 1.18).

We defined the language $\mathrm{L}_{\mathrm{RV}}{ }^{\omega}$ as the multisorted language with

- sorts for $K$ and $\mathrm{RV}_{\delta_{n}}$ for $n<\omega$.
- the ring language for $K$,
- for all $n<\omega$, the (multiplicative) group language as well as the symbol 0 for $\mathrm{RV}_{\delta_{n}}$.
- relation symbols $\oplus_{\delta_{l}, \delta_{m}, \delta_{n}}$ for $n \leq l, m$ integers, function symbols $\mathrm{rv}_{\delta_{n}}: \mathcal{K} \rightarrow \mathrm{RV}_{\delta_{n}}$ and $\operatorname{rv}_{\delta_{n} \rightarrow \delta_{m}}: \operatorname{RV}_{\delta_{n}} \rightarrow \operatorname{RV}_{\delta_{m}}$ for $n>m$.

The structure $\mathcal{K}=\left(K,\left(\operatorname{RV}_{\delta_{n}}\right)_{n<\omega},\left(\oplus_{\delta_{l}, \delta_{m}, \delta_{n}}\right)_{n<l, m},\left(\operatorname{rv}_{\delta_{n}}\right)_{n<\omega},\left(\operatorname{rv}_{\delta_{n} \rightarrow \delta_{m}}\right)_{m<n<\omega}\right)$ becomes a structure in this language where all symbols are interpreted as before. This language is also bi-interpretable (without parameters) with the usual languages of valued fields, e.g. with $\mathrm{L}_{\mathrm{div}}$ (see [17, Proposition 2.8]).

Let $\mathcal{K}$ be any valued field. Let us state few lemmas.
Notation. Let $\delta_{1}, \delta_{2}, \delta_{3} \in \Gamma$ be three values. We write:

$$
\mathrm{WD}_{\delta_{1}, \delta_{2}, \delta_{3}}(\mathbf{x}, \mathbf{y}) \equiv \exists!\mathbf{z} \in \mathrm{RV}_{\delta_{3}} \oplus_{\delta_{1}, \delta_{2}, \delta_{3}}(\mathbf{x}, \mathbf{y}, \mathbf{z})
$$

If the context is clear and in order to simplify notations, we will write:

- $\mathrm{WD}_{\delta_{3}}$ instead of $\mathrm{WD}_{\delta_{1}, \delta_{2}, \delta_{3}}$,
- for any formula $\phi(\mathbf{z})$ with $\mathbf{z} \in \operatorname{RV}_{\delta_{3}}, \mathbf{x} \in \operatorname{RV}_{\delta_{1}}$ and $\mathbf{y} \in \mathrm{RV}_{\delta_{2}}$ :

$$
\phi\left(\operatorname{rv}_{\delta_{3}}(\mathbf{x})+\operatorname{rv}_{\delta_{3}}(\mathbf{y})\right)
$$

or

$$
\phi\left(\mathrm{rv}_{\delta_{3}}(\mathbf{x})+\mathrm{rv}_{\delta_{3}}(\mathbf{y})\right) \wedge \mathrm{WD}_{\delta_{3}}(\mathbf{x}, \mathbf{y})
$$

instead of

$$
\exists \mathbf{z} \in \mathrm{RV}_{\delta_{3}} \oplus_{\delta_{1}, \delta_{2}, \delta_{3}}(\mathbf{x}, \mathbf{y}, \mathbf{z}) \wedge \phi(\mathbf{z}) \wedge \mathrm{WD}_{\delta_{1}, \delta_{2}, \delta_{3}}(\mathbf{x}, \mathbf{y})
$$

Example. Take $K=\mathbb{R}((t))$ the field of power series over the reals endowed with the $t$-adic valuation. Consider $x=t^{2}+t^{3}+t^{4}+t^{5}, x^{\prime}=t^{2}+t^{3}+t^{4}+2 t^{5} \in K, y=-t^{2}-t^{3}+t^{4}-t^{5} \in K$ and $z=2 t^{4}, z^{\prime}=2 t^{4}+t^{5} \in K$.

Then, we have $\operatorname{rv}_{2}(x)=\operatorname{rv}_{2}\left(x^{\prime}\right)$ since $\operatorname{val}\left(x-x^{\prime}\right)=5>\operatorname{val}(x)+2$ but $\mathrm{rv}_{1}(z) \neq \operatorname{rv}_{1}\left(z^{\prime}\right)$ since $\operatorname{val}\left(z-z^{\prime}\right)=5 \ngtr \operatorname{val}(z)+1$. We have:

$$
\vDash \oplus_{2,2,1}\left(\operatorname{rv}_{2}(x), \operatorname{rv}_{2}(y), \operatorname{rv}_{1}(z)\right), \quad \models \oplus_{2,2,1}\left(\operatorname{rv}_{2}(x), \operatorname{rv}_{2}(y), \operatorname{rv}_{1}\left(z^{\prime}\right)\right),
$$

Hence, the sum is not well-defined in $\mathrm{RV}_{1}$ :

$$
\models \neg \mathrm{WD}_{2,2,1}\left(\mathrm{rv}_{2}(x), \mathrm{rv}_{2}(y)\right)
$$

We need to pass to $\mathrm{RV}_{3}$ in order to get a well-defined sum in $\mathrm{RV}_{1}$ :

$$
\begin{gathered}
\models \oplus_{3,3,1}\left(\operatorname{rv}_{3}(x), \operatorname{rv}_{3}(y), \mathrm{rv}_{1}(z)\right), \quad \neg \vDash \oplus_{3,3,1}\left(\mathrm{rv}_{3}(x), \mathrm{rv}_{3}(y), \mathrm{rv}_{1}\left(z^{\prime}\right)\right) \\
\models=\mathrm{WD}_{3,3,1}\left(\operatorname{rv}_{3}(x), \mathrm{rv}_{3}(y)\right)
\end{gathered}
$$

More generally, we have the following proposition:
Proposition 1.48. Let $0 \leq \gamma \leq \delta$ be two elements of $\Gamma_{\geq 0}$ and $\epsilon=\delta-\gamma \geq 0$. Then for every $a, b \in K^{\star}$ :

$$
\mathrm{WD}_{\gamma}\left(\operatorname{rv}_{\delta}(a), \operatorname{rv}_{\delta}(b)\right) \quad \text { if and only if } \quad \operatorname{val}(a+b) \leq \min \{\operatorname{val}(a), \operatorname{val}(b)\}+\epsilon
$$



Proof. Assume $\operatorname{val}(a+b) \leq \min \{\operatorname{val}(a), \operatorname{val}(b)\}+\epsilon$. Let $a^{\prime} \in K$ such that $\operatorname{rv}_{\delta}\left(a^{\prime}\right)=\operatorname{rv}_{\delta}(a)$. This is equivalent to $\operatorname{val}\left(a-a^{\prime}\right)>\operatorname{val}(a)+\delta$, thus we have:

$$
\operatorname{val}\left(a+b-\left(a^{\prime}+b\right)\right)=\operatorname{val}\left(a-a^{\prime}\right)>\operatorname{val}(a)+\delta=\operatorname{val}(a)+\epsilon+\gamma \geq \operatorname{val}(a+b)+\gamma
$$

Hence, $\operatorname{rv}_{\gamma}\left(a^{\prime}+b\right)=\operatorname{rv}_{\gamma}(a+b)$. We have proved the implication from right to left.
Conversely, assume that $\operatorname{val}(a+b)>\min \{\operatorname{val}(a), \operatorname{val}(b)\}+\epsilon$ and $\min \{\operatorname{val}(a), \operatorname{val}(b)\}=\operatorname{val}(a)$. Let $\eta=\operatorname{val}(a+b)+\gamma$ and take any $c \in K$ of valuation $\eta$. Then $\operatorname{rv}_{\delta}(a)=\operatorname{rv}_{\delta}(a+c)$ since $\operatorname{val}(a+c-a)=\eta>\operatorname{val}(a)+\delta$ and $\operatorname{rv}_{\gamma}(a+b) \neq \operatorname{rv}_{\gamma}(a+c+b)$ since $\operatorname{val}(a+c+b-(a+b))=\eta=$ $\operatorname{val}(a+b)+\gamma$.

Remark 1.49. To prove $\operatorname{val}(a+b) \leq \min \{\operatorname{val}(a), \operatorname{val}(b)\}+\epsilon$ with $\epsilon \geq 0$, it is enough to show that $\operatorname{val}(a+b) \leq \operatorname{val}(a)+\epsilon$ (or $\operatorname{val}(a+b) \leq \operatorname{val}(b)+\epsilon)$. Indeed, if $\operatorname{val}(a)=\operatorname{val}(b)$ then this is clear. If $\operatorname{val}(a)<\operatorname{val}(b)$ or $\operatorname{val}(b)<\operatorname{val}(a)$, this is also clear since we have $\operatorname{val}(a+b)=\operatorname{val}(a) \leq \operatorname{val}(a)+\epsilon$ in the first case and $\operatorname{val}(a+b)=\operatorname{val}(b)<\operatorname{val}(a)+\epsilon$ in the second.

The following lemma is immediate:
Lemma 1.50. Let $a, b$ and $c=a-b$ be elements of $K$ and let $\gamma \in \Gamma_{\geq 0}$. At least one of the following holds:

$$
\begin{align*}
& \vDash \mathrm{WD}_{\gamma}\left(\operatorname{rv}_{\gamma}(a), \operatorname{rv}_{\gamma}(b-a)\right)  \tag{1}\\
& \models \mathrm{WD}_{\gamma}\left(\operatorname{rv}_{\gamma}(b), \operatorname{rv}_{\gamma}(a-b)\right) \tag{2}
\end{align*}
$$

Proof. Notice that exactly one of the following occurs:

1. $\operatorname{val} a=\operatorname{val} c<\operatorname{val} b$,
2. $\operatorname{val} b=\operatorname{val} c<\operatorname{val} a$,
3. $\operatorname{val} a=\operatorname{val} b<\operatorname{val} c$,
4. $\operatorname{val} a=\operatorname{val} b=\operatorname{val} c$.


Let $\mathbf{a}=\operatorname{rv}_{\gamma}(a), \mathbf{b}=\operatorname{rv}_{\gamma}(b)$ and $\mathbf{c}=\operatorname{rv}_{\gamma}(c)$. In cases 2,3 and 4, the difference between $\mathbf{a}$ and $\mathbf{c}$ is well-defined.

$$
\vDash \mathrm{WD}_{\gamma}(\mathbf{a},-\mathbf{c})
$$

In cases 1,3 and 4 , the sum of $\mathbf{b}$ and $\mathbf{c}$ is well-defined:

$$
\equiv \mathrm{WD}_{\gamma}(\mathbf{b}, \mathbf{c})
$$

Benign theory of Henselian valued fields Later in this text, we prove transfer principles for some rather nice Henselian valued fields, that we called here 'benign' (see Definition 1.58 below). The goal of this paragraph is to discuss some essential properties that these benign Henselian valued fields share and that we will use for proving Theorem 3.12.

Let $T$ be a (possibly incomplete) theory of Henselian valued fields. We need first to recall the definition of an angular component (or ac-map). It is a group homomorphism usually denoted by ac $:\left(K^{\star}, \cdot\right) \rightarrow\left(k^{\star}, \cdot\right)$ such that $a c_{\mid \mathcal{O}^{\times}}=$res ${ }_{\mathcal{O}} \times$. We also set $\operatorname{ac}(0)=0$. We have the following diagram:


One remarks that an angular component gives a section $a c_{\mathrm{rv}}: \mathrm{RV}^{\star} \rightarrow k^{\star}$ and, as a consequence, the sort $\mathrm{RV}^{\star}$ becomes isomorphic as a group to the direct product $\Gamma \times k^{\star}$. Such a map always exists in an $\aleph_{1}$-saturated valued field $\mathcal{K}$ : as $\mathcal{O}^{\star}$ is a pure subgroup of $K^{\star}$, there is a section $s: \Gamma \rightarrow K^{\star}$ of the valuation (see Fact 1.73). Then, the function ac : $a \mapsto \operatorname{res}(a / s(v(a)))$ is an ac-map. Any theory $T$ of Henselian valued fields in a language $\mathrm{L}_{\Gamma, k}$ admits a natural expansion - denoted by $T_{\mathrm{ac}}$ - in the language $\mathrm{L}_{\Gamma, k, \mathrm{ac}}=\mathrm{L}_{\Gamma, k} \cup\{\mathrm{ac}: K \rightarrow k\}$ by adding the axiom saying that ac is an angular component.

An important model theoretic property is relative quantifier elimination:
$T_{\mathrm{ac}}$ has quantifier elimination (resplendently)
relatively to $\Gamma$ and $k$ in the language $\mathrm{L}_{\Gamma, k, \mathrm{ac}}$.
$T$ has quantifier elimination (resplendently)
relatively to RV in the language $\mathrm{L}_{\mathrm{RV}}$.

Notice that according to the terminology in Paragraph 1.1.1, $\{\Gamma\},\{k\}$ and $\{R V\}$ are closed sets of sorts. Then resplendency automatically follows from relative quantifier elimination (Fact 1.5).

Observation 1.51. $(E Q)_{\mathrm{RV}}$ implies $(E Q)_{\Gamma, k, a c}$.
We include a proof of this observation for completeness.
Proof. We sketch a proof using the usual back-and-forth criterion. We assume (EQ) RV. Consider two models $\mathcal{M}=\left\{K_{M}, \Gamma_{M}, k_{M}\right\}$ and $\mathcal{N}=\left\{K_{N}, \Gamma_{N}, k_{N}\right\}$ of $T$ in the language $\mathrm{L}_{\Gamma, k, \mathrm{ac}}$, and a partial automorphism $f=\left(f_{K}, f_{\Gamma}, f_{k}\right): A=\left(K_{A}, \Gamma_{A}, k_{A}\right) \rightarrow B=\left(K_{B}, \Gamma_{B}, k_{B}\right)$ between a substructure $A \subseteq \mathcal{M}$ and a substructure $B \subseteq \mathcal{N}$. Moreover, we assume $f_{k}$ and $f_{\Gamma}$ to be elementary as morphisms respectively of fields and of ordered abelian groups. We want to extend $f$ to an elementary embedding of $M$ into $N$. By elementarity, we may extend $f_{\Gamma}$ (resp. $f_{k}$ ) to an elementary embedding of ordered abelian groups $\tilde{f}_{\Gamma}: \Gamma_{M} \rightarrow \Gamma_{N}$ (resp. to an elementary embedding of fields $\tilde{f}_{k}: k_{M} \rightarrow k_{N}$ ). Then, by studying quantifier-free formulas, one sees that $\tilde{f}=f \cup \tilde{f}_{\Gamma} \cup \tilde{f}_{k}$ is a partial isomorphism of substructures. Without loss, assume that $\Gamma_{A}=\Gamma_{M}$ and $k_{A}=k_{M}$ and reset the notation. As the ac-map induces a splitting of the exact sequence

$$
1 \rightarrow k^{\star} \xrightarrow{\iota} \mathrm{RV}^{\star} \stackrel{\mathrm{val}_{\mathrm{Fv}}}{\rightarrow} \Gamma \rightarrow 0
$$

we have the bijections $\mathrm{RV}_{M}^{\star} \simeq k_{M}^{\star} \times \Gamma_{M}$ and $\mathrm{RV}_{N}^{\star} \simeq k_{N}^{\star} \times \Gamma_{N}$. Hence, the partial isomorphism $f$ induces an elementary embedding of RV-structure $f_{\mathrm{RV}}:\left(\mathrm{RV}_{M}, \oplus, \cdot, \mathbf{1}, \mathbf{0}\right) \rightarrow\left(\mathrm{RV}_{N}, \oplus, \cdot, \mathbf{1}, \mathbf{0}\right)$, and $f_{K} \cup f_{\mathrm{RV}}$ is a partial isomorphism of substructures in the language $\mathrm{L}_{\mathrm{RV}}$. By relative quantifier elimination down to $\mathrm{RV}, f_{K} \cup f_{\mathrm{RV}}$ extends to an elementary embedding $\tilde{f}=\left(\tilde{f}_{K}, f_{\mathrm{RV}}\right)$ of $\left\{M, \mathrm{RV}_{M}\right\}$ into $\left\{N, \mathrm{RV}_{N}\right\}$. One sees that $\tilde{f}_{K} \cup f_{\Gamma} \cup f_{k}: \mathcal{M} \rightarrow \mathcal{N}$ is an embedding extending the original partial isomorphism $f$. By back-and-forth, $T$ satisfies (EQ) ${ }_{\Gamma, k, \text { ac }}$.

More specifically, we will have to study 1-dimensional definable sets $D \subset K$. Flenner showed in [17] that in Henselian valued fields of characteristic 0, definable sets can be written with field-sorted linear terms (see Fact 1.64). This property will be of essential use. Let us give it also an abbreviation:

Definition 1.52. Let $T$ be the theory of a Henselian valued field $K$ in the language $\mathrm{L}_{\mathrm{RV}}$. We denote by $(\operatorname{Lin})_{\mathrm{RV}}$ the following property: any formula $\phi(x)$ with parameters in $K$ and with $|x|=1$ is equivalent to a formula of the form

$$
\begin{equation*}
\phi_{\mathrm{RV}}\left(\operatorname{rv}\left(x-a_{1}\right), \ldots, \operatorname{rv}\left(x-a_{r}\right), \alpha\right) \tag{3}
\end{equation*}
$$

where $r \in \mathbb{N}$ and $\phi_{\mathrm{RV}}$ is an RV-formula with a tuple of parameters $\alpha \in \operatorname{RV}(K)$ and $a_{1}, \ldots, a_{r} \in K^{1}$.
Notice that it is an improvement of a relative quantifier elimination down to RV for unary-definable sets: the term inside rv is linear in $x$ where (EQ) $)_{\mathrm{RV}}$ gives only a polynomial in $x$.

Its algebraic counterpart seems to be the following:
Definition 1.53. A valued field is called algebraically maximal if it admits no immediate algebraic extension.

In particular, Henselian valued fields of equicharacteristic 0 are algebraically maximal (by the fundamental equality) as well as algebraically closed valued fields. Delon proved that this is actually a first order property. For details, we refer to Delon's thesis [13] and a recent work of Halevi and Hasson in [18].

We show now that algebraically maximal valued fields with quantifier elimination relative to RV enjoy the property $(\mathrm{Lin})_{\mathrm{RV}}$. This fact was suggested by Yatir Halevi. Notice that a similar statement has been proved by Peter Sinclair for valued fields in the Denef-Pas language $\mathrm{L}_{\Gamma, k, a c}$ (see [36, Theorem 2.1.1.]). We thank both of them for their enlightenment.

Let $\mathcal{K}=\left\{(K, \cdot,+, 0,1),(\mathrm{RV}, \mathbf{0}, \mathbf{1}, \cdot, \oplus), \mathrm{rv}: K^{\star} \rightarrow \mathrm{RV}^{\star}\right\}$ be a Henselian valued field viewed as a structure in the language $\mathrm{L}_{\mathrm{RV}}$. Let $\mathbb{K} \succ \mathcal{K}$ be a monster model. Let $x \in \mathbb{K} \backslash K$. We denote by $I_{K}(x)$ the set of values $\{\operatorname{val}(x-a) \mid a \in K\}$. We deduce the following well-known fact from [24, Theorem 1 \& 2]:

Fact 1.54 (Kaplansky). Assume that $\mathcal{K}$ is algebraically maximal. Then $I_{K}(x)$ has no maximum if and only if the extension $K(x) / K$ is immediate.

Lemma 1.55. Assume that $\mathcal{K}$ is algebraically maximal. Let $x \in \mathbb{K} \backslash K$.

- If $I_{K}(x)$ has no maximum, let $\left(\gamma_{i}=\operatorname{val}\left(x-a_{i}\right)\right)_{i \in I}$ be a co-final sequence of values in $I_{K}(x)$. Then the quantifier-free type $\operatorname{qftp}(x / K)$ is implied by the type $\left\{\operatorname{val}\left(x-a_{i}\right)=\gamma_{i}\right\}_{i \in I}$.
- If $I_{K}(x)$ has a maximum, then there is a $\in K$ such that $\operatorname{tp}(\operatorname{rv}(x-a) / \operatorname{RV}(K)$ determined $\mathrm{qftp}(x / K)$. Moreover, $\mathrm{RV}(K(x))$ is generated by $\mathrm{RV}(K)$ and $\operatorname{rv}(x-a)$.

Proof. Assume that $I_{K}(x)$ has no maximum, and let $\left(\gamma_{i}=\operatorname{val}\left(x-a_{i}\right)\right)_{i \in I}$ be a co-final sequence of values in $I_{K}(x)$. Then, the sequence $\left(a_{i}\right)_{i \in I}$ is a pseudo-Cauchy sequence in $K$, with no pseudo-limit in $K$ and which pseudo-converges to $x$. The extension is immediate and $\mathcal{K}$ is algebraically maximal. By [24], the pseudo-Cauchy sequence $\left(a_{i}\right)_{i \in I}$ is of transcendental type. Then, if $x^{\prime} \in \mathbb{K}$ is another pseudo-limit of $\left(a_{i}\right)_{i \in I}$, the two extensions $K(x)$ and $K\left(x^{\prime}\right)$ are isomorphic over $K$. In other words, the quantifier-free type of $x$ over $K$ is uniquely determined by $\left\{\operatorname{val}\left(x-a_{i}\right)=\gamma_{i}\right\}_{i \in I}$.
Assume that $I_{K}(x)$ has a maximum $\gamma=\operatorname{val}(x-a)$. Then, we distinguish two cases.
Case 1: We have that $\operatorname{val}(x-a) \in \Gamma_{K}$.
Then, we have that $\operatorname{rv}(x-a) \notin \operatorname{RV}(K)$, as otherwise, it would exist $b \in K$ such that $\operatorname{val}(x-a-b)>$ $\operatorname{val}(x-a)$, contradicting the maximality of $\operatorname{val}(x-a)$. Let $c \in K$ such that $\operatorname{val}(x-a)=\operatorname{val}(c)$. Then, since $k_{K}$ is relatively algebraically closed in $k_{\mathbb{K}}$ (as $\mathbb{K}$ is an elementary extension of $K$ ) and $\operatorname{since} \operatorname{rv}\left(\frac{x-a}{c}\right)$ is in $k_{\mathbb{K}} \backslash k_{K}$, we have that res $\left(\frac{x-a}{c}\right)$ is transcendental over $k_{K}$. Without loss of generality, assume that $\frac{x-a}{c}=x$. The extension $K(x)$ is the Gauss extension, thus it is unique up to $K$-isomorphism (see e.g. [15]). In particular, the quantifier-free type of $x$ over $K$ is uniquely determined by $\operatorname{tp}(\operatorname{rv}(x) / \operatorname{RV}(K))$. We show now that $\operatorname{RV}(K(x))$ is generated by $\operatorname{RV}(K)$ and $\operatorname{rv}(x)$ in the following sense: Consider $P(X):=\sum_{i<n} a_{i} X^{i}$ a non-trivial polynomial in $K$, and assume that $a_{i_{0}}, \ldots, a_{i_{k-1}}$ are the coefficient of minimum value. Since $\operatorname{rv}\left(x^{n}\right) \in k_{\mathbb{K}}^{\star}$ for all $n$, and since $\operatorname{rv}(x)$ is transcendental over $k_{K}$, we have:

$$
\operatorname{rv}\left(\frac{P(x)}{a_{i_{0}}}\right)=\sum_{j<k} \frac{\operatorname{rv}\left(a_{i_{j}}\right)}{\operatorname{rv}\left(a_{i_{0}}\right)} \operatorname{rv}\left(x^{i_{j}}\right) \in k_{\mathbb{K}}^{\star},
$$

and so

$$
\operatorname{rv}(P(x))=\bigoplus_{j<k} \operatorname{rv}\left(a_{i_{j}}\right) \operatorname{rv}(x)^{i_{j}}=\bigoplus_{i<n} \operatorname{rv}\left(a_{i}\right) \operatorname{rv}(x)^{i} .
$$

Case 2: We have that $\operatorname{val}(x-a) \notin \Gamma_{K}$. Then for all $n \in \mathbb{N}^{\star}, n \cdot \operatorname{val}(x-a) \notin \Gamma_{K}$, as $\mathbb{K}$ is an elementary extension of $K$. Then, for any polynomial $P(x) \in K(x), \operatorname{val}(P(x-a))$ can be expressed in terms of $\operatorname{val}(x-a)$. The isomorphism type of $K(x)$ over $K$ is uniquely determined by $\operatorname{rv}(x-a)=\alpha$, or in other words the quantifier-free type of $x$ over $K$ is uniquely determined by $\operatorname{tp}(\operatorname{rv}(x-a) / \operatorname{RV}(K))$. Without loss of generality, we may assume that $x=x-a$. One sees as well that $\operatorname{RV}(K(x))$ is generated by $\mathrm{RV}(K)$ and $\operatorname{rv}(x)$ : consider $P(X):=\sum_{i<n} a_{i} X^{i}$ a non-trivial polynomial in $K$. As $n \cdot \operatorname{val}(x)$ is not in $\Gamma_{K}$ for all $n$, we have that:

$$
\operatorname{rv}(P(x))=\bigoplus_{i<n} \operatorname{rv}\left(a_{i}\right) \operatorname{rv}\left(x^{i}\right)
$$

Theorem 1.56. Assume that $\mathcal{K}$ is algebraically maximal and admits quantifier elimination relative to RV . Then it also satisfies the property $(\operatorname{Lin})_{\mathrm{RV}}$.

Proof. By compactness, it is enough to show that any (complete) 1-type $p(x)=\operatorname{tp}(b / K)$ over $K$ is determined by formulas of the form 3 .

- If $p(x)=\operatorname{tp}(b / K)$ is a realised type, i.e. $b \in K$, then the type is determined by $\{\operatorname{rv}(x-b)=\mathbf{0}\}$.
- If $K(b) / K$ is immediate, then by the previous lemma, $\operatorname{qftp}(b / K)$ is determined by the type $\left\{\operatorname{val}\left(x-a_{i}\right)=\gamma_{i}\right\}_{i \in I}$, where $\gamma_{i}$ and $a_{i}$ are given by the previous lemma. This can be written in the language $\mathrm{L}_{\mathrm{RV}}$ : for $i \in I$ choose $c_{i} \in K$ of value $\gamma_{i}$. Then

$$
\operatorname{val}\left(x-a_{i}\right)=\gamma_{i} \Leftrightarrow \operatorname{rv}\left(x-a_{i}\right) \oplus \operatorname{rv}\left(c_{i}\right) \neq \operatorname{rv}\left(x-a_{i}\right) \wedge \operatorname{rv}\left(x-a_{i}\right) \oplus \operatorname{rv}\left(c_{i}\right) \neq \operatorname{rv}\left(c_{i}\right)
$$

As $\operatorname{RV}(K(b))=\operatorname{RV}(K)$, and by quantifier elimination relative to the $\operatorname{RV}$-sort, $\left\{\operatorname{val}\left(x-a_{i}\right)=\right.$ $\left.\gamma_{i}\right\}_{i \in I} \cup \operatorname{tp}(\emptyset / \operatorname{RV}(K))$ determines $p(x)=\operatorname{tp}(b / K)$.

- If $K(b) / K$ is non-immediate, then by the previous lemma, there is an $a \in K$ such that $\mathrm{qftp}(b / K)$ is determined by $q(\operatorname{rv}(x-a))$ where $q=\operatorname{tp}(\operatorname{rv}(b-a) / \operatorname{RV})$. As $\mathrm{RV}(K(x))$ is generated by $\operatorname{RV}(K)$ and $\operatorname{rv}(b-a)$, and by quantifier elimination relative to RV, we got that $q(\operatorname{rv}(x-a))$ determines $p(x)=\operatorname{tp}(b / K)$.

Definition 1.57. A valued field of equicharacteristic $p>0$ is said Kaplansky if the value group is $p$-divisible, the residue field is perfect and does not admit any finite separable extensions of degree divisible by $p$.

Definition 1.58. Any $\{\Gamma\}$ - $\{k\}$-enrichment of one of the following theories of Henselian valued fields is called benign:

1. Henselian valued fields of characteristic $(0,0)$,
2. algebraically closed valued fields,
3. algebraically maximal Kaplansky Henselian valued fields.

A model of a benign theory will be called a benign Henselian valued field.
As promised, we have:
Fact 1.59. Benign theories satisfy $(E Q)_{\mathrm{RV}}$ and $(\text { Lin })_{\mathrm{RV}}$.
By the discussion above, they also satisfy $(\mathrm{EQ})_{\Gamma, k, a c}$.
Proof. We just give examples of references for $(\mathrm{EQ})_{\mathrm{RV}}$ and $(\mathrm{EQ})_{\Gamma, k, a c}$. Notice that we might not refer to original proofs. The fact that Henselian valued fields of characteristic $(0,0)$ has property $(\mathrm{EQ})_{\Gamma, k, a c}$ is the classical theorem of Pas. The proof that it has (EQ $)_{\mathrm{RV}}$ is in [17]. Algebraically closed valued fields (in any characteristic) eliminate quantifiers by the theorem of Robinson. One deduces the property $(\mathrm{EQ})_{\Gamma, k, a c}$ from it. One can find a proof that algebraically closed valued fields of any characteristic have $(E Q)_{\mathrm{RV}}$ in [20]. Algebraically maximal Kaplansky valued fields have $(\mathrm{EQ})_{\Gamma, k, a c}$ and $(\mathrm{EQ})_{\mathrm{RV}}$ by [18]. As all these fields are algebraically maximal, they satisfy the condition of Theorem 1.56, and thus enjoy the property $(\mathrm{Lin})_{\mathrm{RV}}$.

Finally, all these properties hold for any $\{\Gamma\}-\{k\}$-enrichment, as it is a particular case of $\{R V\}$ enrichment, and as the sorts $\Gamma, k$ and RV are closed (Fact 1.5).

We will complete our study with some transfer principle for unramified mixed characteristic Henselian valued fields with perfect residue field. As it requires further techniques, it needs to be treated independently. We first need to introduce the Witt vector construction.

Witt vectors In the theory of unramified mixed characteristic Henselian valued fields, we will understand $\mathrm{RV}_{\delta_{n}}$-structures thanks to the well known Witt vector construction. We only introduce some standard notation. For a definition and basic property, one can see [31]. Let $k$ be a field of characteristic p.

Notation (Witt vectors). The ring of Witt vectors over $k$, is denoted by $W(k)$. It admits $k^{\omega}$ as a base set. The residue map $\pi$ is simply the projection to the first coordinate. The natural section of the residue map, the so called Teichmüller lift, is defined as follows:

$$
\begin{aligned}
& \tau: k \longrightarrow W(k) \\
& a \longmapsto[a]:=(a, 0,0, \ldots)^{.}
\end{aligned}
$$

These notions make sense if we restrict the base-set to $k^{n}$. One gets then the truncated ring of Witt vectors of length $n$ denoted by $W_{n}(k)$, as well as its Teichmüller map $\tau_{n}: k \rightarrow W_{n}(k)$.

Observation 1.60. $\left(W_{n}(k),+, \cdot, \pi\right)$ is interpretable in the field $(k,+, \cdot, 1,0)$, with base set $k^{n}$. It is clear that $\operatorname{bdn}\left(W_{n}(k)\right):=\kappa_{\text {inp }}^{1}\left(W_{n}(k)\right) \leq \kappa_{\text {inp }}^{n}(k)$ (we will show that they are in fact equal).

Recall that a $p$-ring is a complete local ring $A$ of maximal ideal $p A$ and perfect residue field $A / p A$. It is strict if $p^{n} \neq 0$ for every $n \in \mathbb{N}$. Here are some basic facts about Witt vectors:

Fact 1.61 (see e.g. [39, Chap. 6]). Recall that the field $k$ is perfect.

1. The ring of Witt vectors $W(k)$ is a strict local $p$-ring of residue field $k$, unique up to isomorphism with these properties.
2. The Teichmüller map is given by the following: let $a \in k$, then $\tau(a)$ is the limit (for the topology given by the maximal ideal $p W(k))$ of any sequence $\left(a_{n}^{p^{n}}\right)_{n<\omega}$ such that $\pi\left(a_{n}\right)^{p^{n}}=a$ for all $n$.
3. In particular $\tau_{n}$ is definable in the structure $\left(W_{n}(k),+, \cdot, \pi\right)$. Indeed $\tau_{n}(a)$ is the (unique) element $a_{n-1}^{p^{n-1}} \in W_{n}$ such that $\pi\left(a_{n-1}\right)^{p^{n-1}}=a$.
4. for $x=\left(x_{n}\right)_{n<\omega} \in W(k)$, one has $x=\sum_{n<\omega}\left[x_{n}^{p^{-n}}\right] p^{n}$.
5. In particular, the map $\chi_{i}: W(k) \rightarrow k, x=\left(x_{0}, x_{1}, \ldots\right) \mapsto x_{i}$ is definable in the structure ( $W(k),+, \cdot, \pi)$. One has indeed

$$
x_{i}=\pi\left(\frac{x-p^{i-1}\left[x_{i-1}^{p^{1-i}}\right]-\cdots-p\left[x_{1}^{p^{-1}}\right]-\left[x_{0}\right]}{p^{i}}\right)^{p^{i}} .
$$

Similarly, for $0 \leq i \leq n-1$ the map $\chi_{n, i}: W_{n}(k) \rightarrow k, x=\left(x_{0}, x_{1}, \ldots, x_{n-1}\right) \mapsto x_{i}$ is definable in $\left(W_{n}(k),+, \cdot, \pi\right)$.

We deduce the following:
Corollary 1.62. - The structure $\left(W_{n}(k),+, \cdot, \pi: W_{n}(k) \rightarrow k\right)$ is bi-interpretable on unary sets with the structure ( $k^{n}, k,+, \cdot, p_{i}, i<n$ ), where $p_{i}: k^{n} \rightarrow k,\left(x_{0}, \ldots, x_{n-1}\right) \mapsto x_{i}$ is the projection map. In other words, there is a bijection $W_{n}(k) \simeq k^{n}$ which leads to identify definable sets.

- Similarly, the structure $(W(k),+, \cdot, \pi: W(k) \rightarrow k)$ is bi-interpretable on unary sets (with no parameters) with the structure $\left(k^{\omega}, k,+, \cdot, p_{i}, i<\omega\right)$ where $p_{i}: k^{\omega} \rightarrow k,\left(x_{0}, x_{1}, \ldots\right) \mapsto x_{i}$.

This corollary, coupled with Fact 1.18, will be one of the main argument to treat our reduction principle in the context of unramified mixed characteristic valued fields with perfect residue field.

Unramified mixed characteristic Henselian valued fields We give a short overview on unramified valued fields, by presenting the similarities with benign valued fields. The (partial) theory of Henselian valued fields of characteristic 0 does not satisfy either $(E Q)_{\Gamma, k, a c}$ or $(E Q)_{\mathrm{RV}}$. We indeed need to get 'information' modulo $\mathfrak{m}_{\delta_{n}}$ in a quantifier-free way. Let us recall the leading term language of finite order:

$$
\mathrm{L}_{\mathrm{RV}}{ }_{<\omega}=\left\{K,\left(\mathrm{RV}_{\delta_{n}}\right)_{n<\omega},\left(\oplus_{\delta_{l}, \delta_{m}, \delta_{n}}\right)_{n<l, m},\left(\operatorname{rv}_{\delta_{n}}\right)_{n<\omega},\left(\mathrm{rv}_{\delta_{n} \rightarrow \delta_{m}}\right)_{m<n<\omega}\right\},
$$

where $\oplus_{n}$ are ternary relation symbols and $\delta_{n}=\operatorname{val}\left(p^{n}\right)$. Let us just define all the analogous properties:

$$
\mathcal{K}, \mathcal{K}^{\prime} \equiv T, \mathcal{K} \subseteq \mathcal{K}^{\prime} \text {, we have } \mathcal{K} \preceq \mathcal{K}^{\prime} \Leftrightarrow \mathrm{RV}_{<\omega} \preceq \mathrm{RV}_{<\omega}^{\prime}
$$

Let us cite two main results in [17]. First we have:
Fact 1.63. [17, Proposition 4.3] Let $T$ be the theory of characteristic 0 Henselian valued fields in the language $\mathrm{L}_{\mathrm{RV}}{ }^{\infty}$. Then $T$ eliminates field-sorted quantifiers.


This result was already proved in [5]. Again, an important consequence is that the multisorted substructure $\left(\left(\operatorname{RV}_{\delta_{n}}\right)_{n<\omega},\left(\oplus_{\delta_{l}, \delta_{m}, \delta_{n}}\right)_{n<l, m},\left(\mathrm{rv}_{\delta_{n} \rightarrow \delta_{m}}\right)_{m<n<\omega}\right)$ is stably embedded and pure. Secondly, we have its one-dimensional improved version:

Fact 1.64. [17, Proposition 5.1] Let $T$ be the theory of Henselian valued fields $K$ of characteristic 0 in the language $\mathrm{L}_{\mathrm{RV}}^{<\omega}$. It has the following property denoted by $(\operatorname{Lin})_{\mathrm{RV}_{<\omega}}$ : any formula $\phi(x)$ with parameters in $K$ and with $|x|=1$ is given by a formula of the form

$$
\begin{equation*}
\phi_{\mathrm{RV}_{\delta_{n}}}\left(\operatorname{rv}\left(x-a_{1}\right), \ldots, \operatorname{rv}\left(x-a_{r}\right), \alpha\right) \tag{4}
\end{equation*}
$$

where $\phi_{\mathrm{RV}_{\delta_{n}}}\left(\mathbf{x}_{1}, \ldots, \mathbf{x}_{r}, \mathbf{y}\right)$ is an $\mathrm{RV}_{\delta_{n}}$-formula, with a tuple of parameters $\alpha \in \mathrm{RV}_{\delta_{n}}(K)$ and $a_{1}, \ldots, a_{r} \in K$ and $r \in \mathbb{N}$.

Again, notice that the improvement comes from the fact that the term inside $\mathrm{rv}_{\delta_{n}}$ is linear in $x$ where Fact 1.63 gives only a polynomial in $x$. These theorems also include the case of equicharacteristic 0 , and it gives the same result as cited above. Indeed, in equicharacteristic 0 , we may identify $\bigcup_{n<\omega} \mathrm{RV}_{\delta_{n}}$ with $R V=R V_{0}$ (Remark 1.46). We continue with a remark on enrichment.

Remark 1.65. Fact 1.63 above holds in any $\mathrm{RV}_{<\omega}$-enrichment of $\mathrm{L}_{\mathrm{RV}}$. Indeed, first note that the $\mathrm{RV}_{<\omega}$-sort is closed in the language $\mathrm{L}_{\mathrm{RV}}^{<\omega}$, i.e. any relation symbol involving a sort $\mathrm{RV}_{\delta_{n}}$ or any function symbol with a domain involving a sort $\mathrm{RV}_{\delta_{n}}$ only involves such sorts. By Fact 1.5, the theory $T$ of Henselian valued fields of characteristic 0 also eliminates quantifiers resplendently relative to $\mathrm{RV}_{<\omega}$. In other words, given an $\mathrm{RV}_{<\omega \text {-enrichment }} \mathrm{L}_{\mathrm{RV}}^{<\omega, e}$, any complete $\mathrm{L}_{\mathrm{RV}}{ }_{<\omega, \mathrm{e}}$-theory $T_{e} \supset T$ eliminates field-sorted quantifiers. A careful reading of Flenner's proof give us that Fact 1.64 holds resplendently as well.

Now, let us discuss more specifically on the unramified mixed characteristic cases. We denote by $T$ the theory of unramified mixed characteristic Henselian valued fields with perfect residue field. We assume now that $\mathcal{K}$ is such a valued field. There is by definition a smallest positive value $\operatorname{val}(p)=\delta_{1}$, that we denote by 1 .

Notation: Notice that $\mathfrak{m}=p \mathcal{O}$ and in general that $\mathfrak{m}_{\delta_{n}}=\mathfrak{m}^{n+1}=p^{n+1} \mathcal{O}$ for all $n \geq 0$. We will write $\mathfrak{m}^{n+1}$ instead of $\mathfrak{m}_{\delta_{n}}, \mathcal{O}_{n+1}$ instead of $\mathcal{O}_{\delta_{n}}$ and $\mathrm{RV}_{n+1}$ instead of $\mathrm{RV}_{\delta_{n}}=K^{\star} / 1+p^{n+1} \mathcal{O}$. The projection map $\operatorname{res}_{\delta_{n}}: \mathcal{O} \rightarrow \mathcal{O}_{\delta_{n}}$ is written $\operatorname{res}_{n+1}: \mathcal{O} \rightarrow \mathcal{O}_{n+1}$ etc. The idea is to denote by $\mathrm{RV}_{n}$ the $n^{\text {th }}$ RV-sort, as this makes sense in unramified (or finitely ramified) mixed characteristic valued fields. The purpose is also to fit with the usual notation, and it will help to simplify the notation, although this convention contradicts the previous one ( $R V_{0}$ where 0 stands for the value $0 \in \Gamma$ is now $R V_{1}$, the first RV-sort).

In this context, let us define the angular component of degree $n$ :
Definition 1.66. Let $n$ be an integer greater than 0 . An angular component of order $n$ is a homomorphism ac ${ }_{n}: K^{\star} \rightarrow \mathcal{O}_{n}^{\times}$such that for all $u \in \mathcal{O}^{\times}, \operatorname{ac}_{n}(u)=\operatorname{res}_{n}(u)$. A system of angular component maps $\left(a c_{n}\right)_{n<\omega}$ is said to be compatible if for all $n$, $\pi_{n} \circ \mathrm{ac}_{n+1}=\mathrm{ac}_{n}$ where $\pi_{n}: \mathcal{O}_{n+1} \rightarrow \mathcal{O}_{n} \simeq \mathcal{O}_{n+1} / p^{n} \mathcal{O}_{n+1}$ is the natural projection.


The convention is to contract $\mathrm{ac}_{1}$ to ac and $\pi_{1}$ to $\pi$. Then, let us complete the diagram given in Paragraph 1.2.1:


A section $s: \Gamma \rightarrow K^{\star}$ of the valuation gives immediately a compatible system of angular components (defined as ac $\left.{ }_{n}:=a \in K^{\star} \mapsto \operatorname{res}_{n}(a / s(v(a)))\right)$. As $\mathcal{O}^{\times}$is a pure subgroup of $K^{\star}$, such a section exists when $\mathcal{K}$ is $\aleph_{1}$-saturated (see Fact 1.73). As always, we assume that $\mathcal{K}$ is sufficiently saturated and we fix a compatible sequence $\left(\mathrm{ac}_{n}\right)_{n}$ of angular components.

We denote by $T_{a c<\omega}$ the extension of $T$ to the language $\mathrm{L}_{\Gamma, k, \mathrm{ac}<\omega}:=\mathrm{L}_{\Gamma, k} \cup\left\{\mathcal{O}_{n}, \mathrm{ac}_{n}: K \rightarrow \mathcal{O}_{n}, n \in\right.$ $\mathbb{N}\}$ where $\mathrm{ac}_{n}$ are interpreted as compatible angular components of degree $n$ (see for instance [2]).

The following proposition is well known and has been used for example in [6, Corollary 5.2]. It states how the structure $\mathrm{RV}_{n}$ and the truncated Witt vectors $W_{n}$ are related.

Proposition 1.67. 1. The residue ring $\mathcal{O}_{n}$ of order $n$ is isomorphic to $W_{n}(k)$, the set of truncated Witt vectors of length $n$.
2. The kernel of the valuation val : $\mathrm{RV}_{n}^{\star} \rightarrow \Gamma$ is given by $\mathcal{O}^{\times} /\left(1+\mathfrak{m}^{n}\right) \simeq\left(\mathcal{O} / \mathfrak{m}^{n}\right)^{\times}$. It is isomorphic to $W_{n}(k)^{\times}$, the set of invertible elements of $W_{n}(k)$.
Proof. It is clear that (2) follows from (1) as $\mathcal{O}^{\times} /\left(1+\mathfrak{m}^{n}\right) \simeq\left(\mathcal{O} / \mathfrak{m}^{n}\right)^{\times}$.
Now, we prove (1) for any discrete value group $\Gamma$. Consider

$$
W^{\prime}:=\lim _{n<\omega} \mathcal{O}_{n} \subset \prod_{n<\omega} \mathcal{O}_{n}
$$

the inverse limit of the $\mathcal{O}_{n}$ 's. It is:

- strict, i.e. $p^{n} \neq 0$ in $W^{\prime}$ for every $n<\omega$, as $\pi_{n+1}\left(p^{n}\right) \neq 0$ in $\mathcal{O}_{n+1}=\mathcal{O} / p^{n+1} \mathcal{O}$,
- local, as a projective limit of the local rings $\mathcal{O}_{n}$,
- a $p$-ring. Its maximal ideal is $p W^{\prime}$, it is complete as projective limit, and its residue field is the perfect field $k$.

By uniqueness, $W^{\prime}$ is isomorphic to $W(k)$, the ring of Witt vectors over $k$. One just has to notice that $W^{\prime} / p^{n} W^{\prime} \simeq \mathcal{O} / p^{n} \mathcal{O}$ and it follows easily that $\mathcal{O}_{n} \simeq W_{n}(k)$ for every $n<\omega$.

Note. In the above proof, one can also recover $W(k)$ by considering the coarsening $\dot{K}$ of $K$ by the convex subgroup $\mathbb{Z} \cdot 1$. Indeed, if we denote by $K^{\circ}$ the residue field of the coarsening, as $\mathcal{K}$ is saturated enough, one has $\lim _{n<\omega} \mathcal{O}_{n} \simeq \mathcal{O}\left(K^{\circ}\right)($ see [5] $)$.

Fact 1.68 (Bélair [6]). The theory $T_{a c_{<\omega}}$ of Henselian mixed characteristic valued fields with perfect residue field and with angular components eliminates field-sorted quantifiers in the language $\mathrm{L}_{\Gamma, k, \mathrm{ac}<\omega}{ }^{4}$

Notice that in [6], Bélair doesn't assume that the residue field $k$ is perfect, but it is indeed necessary in order to identify the ring $\mathcal{O}_{n}:=\mathcal{O} / \mathfrak{m}^{n}$ with the truncated Witt vectors $W_{n}(k)$. This implies as well that the residue field $k$ and the value group $\Gamma$ are pure sorts, and are orthogonal. This can be seen by analysing field-sorted-quantifier-free formulas, and by noticing that $\mathcal{O}_{n} \simeq W_{n}(k)$ is interpretable in $k$ (Corollary 1.62).

By analogy with the previous paragraph, we name the following properties:

$$
T_{a c_{<\omega}} \text { eliminates } K \text {-sorted quantifiers in the language } \mathrm{L}_{\mathrm{ac}<\omega}
$$

$T$ has quantifier elimination (resplendently) relatively to $\mathrm{RV}_{<\omega}$.

$$
(\mathrm{EQ})_{\mathrm{RV}_{<\omega}}
$$

Again, notice that $\mathrm{RV}_{<\omega}=\bigcup_{n<\omega} \mathrm{RV}_{n}$ is a closed set of sorts.
To sum up, we have:
Fact 1.69. The theory of unramified mixed characteristic Henselian valued fields with perfect residue field satisfies $(E Q)_{\mathrm{RV}_{<\omega}},(E Q)_{\Gamma, k, a c_{<\omega}},(\text { Lin })_{\mathrm{RV}_{<\omega}}$.

### 1.2.2 Abelian groups

We conclude these preliminaries with some facts on abelian groups. They will be used in Section 2 . We are specifically interested in abelian groups for one main reason: we have to understand the structure of pure short exact sequences of abelian groups in order to produce our reduction principles for benign Henselian valued fields. As we obtain also reduction principles for such short exact sequences, we will take the occasion to apply it on explicit examples.

Notation. We recall some standard notation:

- $\mathbb{Z}\left(p^{n}\right)$ is the cyclic group of $p^{n}$ elements,
- $\mathbb{Z}_{(p)}$ is the additive group of the integers localised in $(p)$,
- $\mathbb{Z}\left(p^{\infty}\right)$ is the Prüfer $p$-group, .

[^4]For an abelian group $A$, we denote by $A^{(\omega)}$ the direct product of $\omega$ copies of $A$.
We have:
Fact 1.70 ( $[29$, Theorem $2 \mathbb{Z} 1])$. Let $A$ be an abelian group. Let $P_{n}$ be a predicate for $n$-divisibility in A. Then $\left\{A,+,-, 0, P_{n}, n \in \mathbb{N}_{>1}\right\}$ eliminates quantifiers.

The burden (or equivalently by stability, the dp-rank) of pure abelian groups has been computed in terms of their Szmielew invariants by Halevi and Palacín in [19]. We borrow from their work the following proposition, which says that a useful criterion to witness inp-patterns is a characterisation in the case of one-based groups (and in particular, in the case of unenriched abelian groups):

Proposition 1.71 ([19, Proposition 3.4]). A stable one-based group admits an inp-pattern of depth $\kappa$ if and only if there exists $\operatorname{acl}^{e q}(\emptyset)$-definable subgroups $\left(H_{\alpha}\right)_{\alpha<\kappa}$ such that for any $i_{0}<\kappa$, one has:

$$
\left[\bigcap_{\alpha \neq i_{0}} H_{\alpha}: \bigcap_{\alpha} H_{\alpha}\right]=\infty
$$

If $\left(b_{\alpha, j}\right)_{j<\omega}$ are representatives of pairwise distinct classes of $\bigcap_{\alpha \neq i_{0}} H_{\alpha}$ modulo $\bigcap_{\alpha} H_{\alpha}$, an inp-pattern of depth $\kappa$ is given by $\left\{x \in b_{\alpha, j} H_{\alpha}\right\}_{\alpha<\kappa, j<\omega}$.

We will use this criterion to provide examples to Theorem 2.2.

## Quantifier elimination result in pure short exact sequences

Definition 1.72. Let $B$ a group and $A$ a subgroup. We say that $A$ is a pure subgroup of $B$ if for all $a$ in $A, n \in \mathbb{N}, a$ is $n$-divisible in $B$ if and only if $a$ is $n$-divisible in $A$.

We recall the following fundamental fact:
Fact 1.73. Let $\mathcal{M}$ be an $\aleph_{1}$-saturated structure, and let $A, B$ be two definable abelian groups, and assume that $A$ is a pure subgroup of $B$. Then the exact sequence of abelian groups $0 \rightarrow A \rightarrow B \rightarrow$ $B / A \rightarrow 0$ splits: there is a group homomorphism $\alpha: B \rightarrow A$ such that $\alpha_{\mid A}$ is the identity on $A$. In such case, $B$ is isomorphic as a group to $A \times B / A$.

More precisely, it is an immediate corollary of a more general statement on pure-injectivity. See [7, Theorem 20 p.171].

Assume that we have a pure short exact sequence of abelian groups

$$
0 \longrightarrow A \xrightarrow{\iota} B \xrightarrow{\nu} C \longrightarrow 0 \text {. }
$$

(meaning that $\iota(A)$ is a pure subgroup of $B$ ). We treat it as a three-sorted structure $(A, B, C, \iota, \nu)$, with a group structure for all sorts. In fact, in our main applications, we will consider such a sequence with more structure on $A$ and $C$. Let us explicitly state all results resplendently, by working in an enriched language. So, let $\mathcal{M}=(A, B, C, \iota, \nu, \ldots)$ be an $\{A\}$-enrichment of a $\{C\}$-enrichment (for short: an $\{A\}$ - $\{C\}$-enrichment) of the exact sequence in a language that we will denote by $L$, and we denote its theory by $T$. We will always assume that $\mathcal{M}$ is sufficiently saturated ( $\aleph_{1}$ saturated will be enough).
Hypothesis of purity implies the exactness of the following sequences for $n \in \mathbb{N}$ :

$$
0 \longrightarrow A / n A \xrightarrow{\iota_{n}} B / n B \xrightarrow{\nu_{n}} C / n C \longrightarrow 0
$$

One has indeed that

$$
\frac{A+n B}{n B} \simeq \frac{A}{A \cap n B}=\frac{A}{n A}
$$

We consider for $n \geq 0$ the following maps:

- the natural projections $\pi_{n}: A \rightarrow A / n A$,
- the map

$$
\begin{aligned}
\rho_{n}: B & \rightarrow \\
b & \mapsto\left\{\begin{array}{c}
A / n A \\
0_{A / n A} \text { if } b \notin \nu^{-1}(n C) \\
\iota_{n}^{-1}(b+n B) \text { otherwise }
\end{array}\right.
\end{aligned}
$$

where $0_{A / n A}$ is the zero element of $A / n A$ (often denoted by 0 ). Then let us consider the language

$$
\mathrm{L}_{q}=\mathrm{L} \cup\left\{A / n A, \pi_{n}, \rho_{n}\right\}_{n \geq 0}
$$

and let $T_{q}$ be the natural extension of the theory $T$. By $<A>$, we denote the set of sorts containing $A, A / n A$ and the new sorts possibly coming from the $A$-enrichment. Similarly, let $<C>$ be the set of sorts containing $C$ and the new sorts possibly coming from the $C$-enrichment. By $A$ or $A$-sort and $C$ or $C$-sort, we will abusively refer to $<A>$ and $<C>$ respectively, and similarly for $A$-formulas and $C$-formulas. Aschenbrenner, Chernikov, Gehret and Ziegler prove the following result:

Fact 1.74 ([3, Theorem 4.2]). The theory $T_{q}$ (resplendently) eliminates B-sorted quantifiers.


More precisely, all $\mathrm{L}_{q}$-formulas $\phi(x)$ with a tuple of variables $x \in B^{|x|}$ are equivalent to boolean combinations of formulas of the form:

1. $\phi_{C}\left(\nu\left(t_{0}(x)\right), \ldots, \nu\left(t_{s-1}(x)\right)\right)$ where $t_{i}(x)$ 's are terms in the group language, and $\phi_{C}$ is a $C$ formula,
2. $\phi_{A}\left(\rho_{n_{0}}\left(t_{0}(x)\right), \ldots, \rho_{n_{s-1}}\left(t_{s-1}(x)\right)\right)$ where the $t_{i}(x)$ 's are terms in the group language, where $s, n_{0}, n_{1}, \ldots, n_{s-1} \in \mathbb{N}$, and where $\phi_{A}$ is an $A$-formula.
In particular there is no occurrence of the symbol $\iota$.
In particular, notice that the formula $t(x)=0$ is equivalent to $\nu(t(x))=0 \wedge \rho_{0}(t(x))=0$ and $\exists y n y=t(x)$ is equivalent to $\exists y_{C} n y_{C}=\nu(t(x)) \wedge \rho_{n}(t(x))=0$.

We have:
Corollary 1.75. In the theory $T_{q},<A>$ and $<C>$ are stably embedded, pure (see Definition 1.7) and orthogonal to each other.

In fact, it can be easily deduced from the existence of a section. The following proof is more technical but highlights the fact that one does not need the function $\iota$ in order to express definable sets in $\bigcup_{n<\omega} A / n A$.

Proof. In this proof, $A$ (resp. $C$ ) abusively refer to the union of the sorts $<A>$ (resp. $<C>$ ). The $C$-sort is pure and stably embedded by Fact 1.74 and closedness of $C$. It is also clear for the sort $A$, even if $A$ is not a closed sort ${ }^{5}$ : one only needs to deal with the map $\iota: A \rightarrow B$. If $D$ is a definable set in $A^{\left|x_{A}\right|}$, it is given by a disjunction of formulas of the form

$$
\begin{aligned}
\phi\left(x_{A}\right)= & \phi_{A}\left(\rho_{n_{0}}\left(k_{0} \iota\left(t_{0}\left(x_{A}\right)\right)+b_{0}\right), \ldots, \rho_{n_{s-1}}\left(k_{s-1} \iota\left(t_{s-1}\left(x_{A}\right)\right)+b_{s-1}\right), a\right) \\
& \wedge \phi_{C}\left(\nu\left(\iota\left(t_{0}\left(x_{A}\right)\right)\right), \ldots, \nu\left(\iota\left(t_{s-1}\left(x_{A}\right)\right)\right), c\right) .
\end{aligned}
$$

[^5]where $x_{A}$ is a tuple of $A$-variables, the $t_{i}\left(x_{A}\right)$ 's are terms in the group language, $s, k_{0}, \ldots, k_{s-1} \in \mathbb{N}$, $b_{0}, \ldots, b_{s-1} \in B$, and $a \in A$ and $c \in C$ are tuples of parameters (notice that we also used that $\iota$ and $\nu$ are morphisms). We apply now the following transformation in order to get a new formula $\phi^{\prime}\left(x_{A}\right)$ :

- For $l<s$, if $\nu\left(b_{l}\right) \notin n_{l} C$, then replace $\rho_{n}\left(k_{l} \iota\left(t_{l}\left(x_{A}\right)\right)+b_{l}\right)$ by $0_{A / n_{l} A}$.
- For $l<s$, if $\nu\left(b_{l}\right) \in n_{l} C$, replace $\rho_{n_{l}}\left(k_{l} \iota\left(t_{l}\left(x_{A}\right)\right)+b_{l}\right)$ by $k_{l} \pi_{n_{l}}\left(t_{l}\left(x_{A}\right)\right)+\rho_{n_{l}}\left(b_{l}\right)$.
- Replace $\nu\left(\iota\left(t_{l}\left(x_{A}\right)\right)\right.$ by $0_{C}$.

We obtain a pure $A$-formula $\phi^{\prime}\left(x_{A}\right)$ such that $\phi^{\prime}\left(A^{\left|x_{A}\right|}\right)=\phi\left(A^{\left|x_{A}\right|}\right)$. Orthogonality can also be proved similarly.

## 2 Burden of pure short exact sequences of abelian groups

We prove in this section that the burden of a pure short exact sequence of abelian groups

$$
0 \rightarrow A \rightarrow B \rightarrow C \rightarrow 0
$$

can be computed in $A$ and $C$ (Theorem 2.2). This result is motivated by valued fields, as an RVstructure can be seen as an enrichment of such. It is one of the main element for our proof of Theorem 3.12 and Theorem 3.21.

### 2.1 Reduction

As in the paragraph 1.2.2, we consider a pure exact sequence $\mathcal{M}$ of abelian groups

$$
0 \longrightarrow A \xrightarrow{\iota} B \xrightarrow{\nu} C \longrightarrow 0,
$$

in an $\{A\}$ - $\{C\}$-enriched language L . In the following paragraph, we compute the burden of the structure $\mathcal{M}$ in terms of burden of $A$ and that of $C$ (in their induced structure). By bi-interpretability on unary sets, one can also consider it as a one-sorted structure $A \subset B$ where $A$ is given by a predicate. It follows indeed from Fact1.18 that $\operatorname{bdn}(A \rightarrow B \rightarrow C)=\operatorname{bdn}(B, A)$. We often prefer the point of view of an exact sequence as it is more relevant for the computation of the burden. We write indifferently $\operatorname{bdn}(\mathcal{M}), \operatorname{bdn}(A \rightarrow B \rightarrow C)$ or $\operatorname{bdn}(B)$, as the sort $B$ is understood as a sort of $\mathcal{M}$ with its full induced structure.

Notice that in the case where $B / n B$ and $B_{[n]}$ are finite for all $n$ and $C$ is torsion free, a straight forward generalisation of [11, Proposition 4.1] gives that $\operatorname{bdn}(A \rightarrow B \rightarrow C)=\max (\operatorname{bdn}(A), \operatorname{bdn}(C))$. We will see that one can get rid of these hypothesis and obtain a more general result using Fact 1.74. We first show a trivial bound:

Fact 2.1 (Trivial bound). Assume there is a section of the group morphism $\nu: B \rightarrow C$. Consider $\mathrm{L}_{s}$ the language L augmented by a symbol s, and interpret it by this section of $\nu$.

$$
0 \longrightarrow A \xrightarrow{\iota} B \stackrel{\stackrel{s}{\nu}}{\longrightarrow} C \longrightarrow
$$

We have $\operatorname{bdn}_{\mathrm{L}}(C)=\operatorname{bdn}_{\mathrm{L}_{s}}(C)$ and $\operatorname{bdn}_{\mathrm{L}}(A)=\operatorname{bdn}_{\mathrm{L}_{s}}(A)$ as well as the following:

$$
\max \left\{\operatorname{bdn}_{\mathrm{L}}(A), \operatorname{bdn}_{\mathrm{L}}(C)\right\} \leq \operatorname{bdn}_{\mathrm{L}}(B) \leq \operatorname{bdn}_{\mathrm{L}_{s}}(B)=\operatorname{bdn}_{\mathrm{L}}(A)+\operatorname{bdn}_{\mathrm{L}}(C)
$$

Proof. The two first equalities are clear since $A$ and $C$ are stably embedded (and orthogonal) in both languages. The inequality $\max \left\{\operatorname{bdn}_{\mathrm{L}}(A), \operatorname{bdn}_{\mathrm{L}}(C)\right\} \leq \operatorname{bdn}_{\mathrm{L}}(B)$ is obvious. As the burden only grows when we add structure, the inequality $\operatorname{bdn}_{\mathrm{L}}(B) \leq \operatorname{bdn}_{\mathrm{L}_{s}}(B)$ is also clear. The last equality come from the fact that in the language $\mathrm{L}_{s}$, the structure $A \rightarrow B \rightarrow C$ and the structure $\left\{A \times C, A, C, \pi_{A}\right.$ : $\left.A \times C \rightarrow A, \pi_{C}: A \times C \rightarrow C\right\}$ are bi-interpretable on unary sets. We conclude by Proposition 1.24 and Fact 1.18.

Theorem 2.2. Consider an $\{A\}-\{C\}$-enrichment of a pure exact sequence $\mathcal{M}$ of abelian groups

$$
0 \longrightarrow A \xrightarrow{\iota} B \xrightarrow{\nu} C \longrightarrow 0,
$$

in a language L. Let $\mathcal{D}=\mathcal{D}(x)$ be the set of formulas in the pure language of groups which are conjunction of formulas of the form $\exists y n x=m y$ for $n, m \in \mathbb{N}$. For $D(x) \in \mathcal{D}$ and $A$ an abelian group, $D(A)$ is an subgroup of $A$, and we have

$$
\operatorname{bdn} \mathcal{M}=\max _{D \in \mathcal{D}}(\operatorname{bdn}(A / D(A))+\operatorname{bdn}(D(C)))
$$

In particular:

- If $A / n A$ is finite for all $n \geq 1$, then

$$
\operatorname{bdn}(\mathcal{M})=\max _{k \in \mathbb{N}}\left(\operatorname{bdn}(k A)+\operatorname{bdn}\left(C_{[k]}\right)\right)
$$

where $C_{[k]}:=\{c \in C \mid k c=0\}$ is the subgroup of $k$-torsion.

- If $C$ has finite $k$-torsion of all $k \geq 1$, then

$$
\operatorname{bdn}(\mathcal{M})=\max _{n \in \mathbb{N}}(\operatorname{bdn}(A / n A)+\operatorname{bdn}(n C))
$$

- If $C$ has finite $n$-torsion and $A / n A$ is finite for all $n \geq 1$, then

$$
\operatorname{bdn} \mathcal{M}=\max (\operatorname{bdn}(A), \operatorname{bdn}(C))
$$

To clarify, for $D(x) \in \mathcal{D}, \operatorname{bdn}(A / D(A))$ can be computed in $A / D(A)$ endowed with its induced structure, i.e. it is the supremum of depth of patterns $P\left(x_{D}\right)$ with $x_{D}$ an $A / D(A)$ variable within the structure $\left\{(A, 0,+, \cdots),(A / D(A), 0,+), \pi_{D}: A \rightarrow A / D(A)\right\}$, where $\cdots$ denotes the enriched structure in $A$.

Remark 2.3. - This result is also valid with the Adler's convention (see Definition 1.31):

$$
\operatorname{bdn}^{\star}(\mathcal{M})=\max _{D \in \mathcal{D}}\left(\operatorname{bdn}^{\star}(A / D(A))+\operatorname{bdn}^{\star}(D(C))\right)
$$

- If $\operatorname{bdn}(A)$ or $\operatorname{bdn}(C)$ is infinite (or equals $\aleph_{0-}$ ), then this is simply the trivial bound in Fact 2.1 (as then $\max (\operatorname{bdn}(A), \operatorname{bdn}(C))=\operatorname{bdn}(A)+\operatorname{bdn}(C))$. Recall that a section exists in a $\aleph_{1}$-saturated model as $A$ is pure in $B$ (Fact 1.73).
- The maximum is always attained by at least one $D \in \mathcal{D}$ : if $\operatorname{bdn}(A)$ and $\operatorname{bdn}(C)$ are finite, this is trivial. If $\operatorname{bdn}(A)$ or $\operatorname{bdn}(C)$ is infinite, then $\operatorname{bdn}(A)$ or $\operatorname{bdn}(C)$ (corresponding terms for resp. $D \equiv x=0$ and $D \equiv x=x$ ) realises the maximum by the previous point.
- By Fact 1.26, if $C$ has finite $n$-torsion for $n \in \mathbb{N}$, one has $\operatorname{bdn}(n C)=\operatorname{bdn}(C)$.
- If $m \mid n$, then $A / m A$ can be seen as a quotient of $A / n A$ and naturally, one has $\operatorname{bdn}(A / m A) \leq$ $\operatorname{bdn}(A / n A)$.

In the case that the sequence is unenriched, this can gives us absolute results: assume that the induced structures on $A / D(A)$ and $D(C)$ are the structures of groups for every $D \in \mathcal{D}$. Then, Proposition [19, Theorem 1.1.] together with Theorem 2.2 gives us a computation of $\operatorname{bdn}(A \rightarrow B \rightarrow C)$ in term of Szmielew invariants of $A$ and $C$. We don't attempt to write a closed formula. Nonetheless, here are some examples:
Examples. We consider the following pairs of abelian groups $(A \subset B)$ with quotient $C$ :

- $B=\mathbb{Z}_{(2)}^{(\omega)} \oplus \mathbb{Z}_{(3)}^{(\omega)}$ and $A=\mathbb{Z}_{(2)}^{(\omega)} \oplus\{0\}$. One has $\operatorname{bdn}(A / n A)+\operatorname{bdn}(n C)=0+1=1$ for all $2 \nmid n$, and $\operatorname{bdn}(A / 2 n A)+\operatorname{bdn}(2 n C)=1+1=2$, which leads to $\operatorname{bdn} \mathcal{M}=2$. This can already be deduced from Halevi and Palacín's work: the sort $B$, equipped only with its group structure, is already of burden 2. By the trivial bound, the structure $\mathcal{M}$ is also of burden 2 .
- $B=\mathbb{Z}_{(2)}^{(\omega)} \oplus \mathbb{Z}(2)^{(\omega)}$ and $A=\mathbb{Z}_{(2)}^{(\omega)} \oplus 0$. Then $\operatorname{bdn}(\mathcal{M})=2$ (take $D(C)=C_{[2]}$ in Theorem 2.2). Again, by [19], the burden of $B$ as a pure abelian group is already 2 .
- $B=\mathbb{Z}_{(2)}^{(\omega)} \oplus \mathbb{Z}_{(3)}^{(\omega)} \oplus \mathbb{Z}_{(2)}^{(\omega)} \oplus \mathbb{Z}_{(3)}^{(\omega)}$ and $A=\mathbb{Z}_{(2)}^{(\omega)} \oplus \mathbb{Z}_{(3)}^{(\omega)} \oplus\{0\} \oplus\{0\}$. Then $\operatorname{bdn}(\mathcal{M})=4$ (take $D(C)=6 C$ in Theorem 2.2). In term of subgroups, one can consider the subgroups $A+4 B$, $A+9 B, 2 B$ and $3 B$. The intersection is

$$
2 \mathbb{Z}_{(2)}^{(\omega)} \oplus 3 \mathbb{Z}_{(3)}^{(\omega)} \oplus 4 \mathbb{Z}_{(2)}^{(\omega)} \oplus 9 \mathbb{Z}_{(3)}^{(\omega)} .
$$

One may see that these groups satisfy Proposition 1.71.

- $A=\mathbb{Z}\left(2^{\infty}\right), C=\mathbb{Z}_{(2)}^{(\omega)} \oplus \mathbb{Z}_{(3)}^{(\omega)}$ and $B=A \times C$. One can see that $\operatorname{bdn}(C \rightarrow B \rightarrow A)=\operatorname{bdn}(C, B)=$ 3. This equality is witnessed by the subgroups $2 B, 3 B$ and $C+B_{[2]}$. However, by Theorem 2.2, $\operatorname{bdn}(A \rightarrow B \rightarrow C)=\operatorname{bdn}(C)=2$ as $A / n A=\{0\}$ for all $n \geq 1$ and $C_{[k]}=0$ for all $k \geq 1$.

Proof of Theorem 2.2. By Fact 1.18, we can work in the language $\mathrm{L}_{q}$ and use Fact 1.74. Recall that we abusively refer the union of sorts $\{A / n A\}_{n \in N}$ as the sort $A$.

The purity of the exact sequence

$$
0 \longrightarrow A \xrightarrow{\iota} B \xrightarrow{\nu} C \longrightarrow 0,
$$

implies for all $n$ the following exact sequences:

$$
\begin{gathered}
0 \longrightarrow A / n A \xrightarrow{\iota_{n}} B / n B \xrightarrow{\nu_{n}} C / n C \longrightarrow 0, \\
0 \longrightarrow A_{[n]} \xrightarrow{\iota^{[n]}} B_{[n]} \xrightarrow{\nu^{[n]}} C_{[n]} \longrightarrow 0,
\end{gathered}
$$

and

$$
0 \longrightarrow A / A_{[n]} \xrightarrow{\iota^{n}} B / B_{[n]} \xrightarrow{\nu^{n}} C / C_{[n]} \longrightarrow 0 \text {. }
$$

All these sequences are again pure, and we can keep going. We get more generally the following:
Fact 2.4. For all $D \in \mathcal{D}$ we have the following pure exact sequences:

$$
0 \longrightarrow D(A) \xrightarrow{\iota^{D}} D(B) \xrightarrow{\nu^{D}} D(C) \longrightarrow 0,
$$

and

$$
0 \longrightarrow A / D(A) \xrightarrow{\iota_{D}} B / D(B) \xrightarrow{\nu_{D}} C / D(C) \longrightarrow 0 .
$$

Proof. By Fact 1.73 , an $\aleph_{1}$-saturated extension $0 \rightarrow A^{\prime} \rightarrow B^{\prime} \rightarrow C^{\prime} \rightarrow 0$ splits. Then we clearly have the exactness of the following sequence:

$$
0 \longrightarrow D\left(A^{\prime}\right) \longrightarrow D\left(B^{\prime}\right) \longrightarrow D\left(C^{\prime}\right) \longrightarrow 0
$$

and

$$
0 \longrightarrow A^{\prime} / D\left(A^{\prime}\right) \longrightarrow B^{\prime} / D\left(B^{\prime}\right) \longrightarrow C^{\prime} / D\left(C^{\prime}\right) \longrightarrow 0
$$

As they are first order properties in the language $\mathcal{L}$, we deduce the fact.
It follows that if $\nu(b) \in D(C)$, there is a unique $a+D(A) \in A / D(A)$ such that $i(a)+D(B)=$ $b+D(B)$. We denote by $\rho_{D}$ the following map:

$$
\begin{aligned}
\rho_{D}: B & \rightarrow \\
b & \mapsto \begin{cases}\iota_{D}{ }^{-1}(b+D(B)) & \text { if } \nu(b) \in D(C) \\
0 & \text { otherwise }\end{cases}
\end{aligned}
$$

Notice that it is interpretable in the language $\mathcal{L}$.
By Lemma 1.26, if $A$ is finite, we get that $\operatorname{bdn}(B)=\operatorname{bdn}(C)$ and of course $\operatorname{bdn}(A / D(A))=0$ for all $D \in \mathcal{D}$. Assume that $A$ is infinite. As $A$ and $C$ are orthogonal, so are in particular $A / D(A)$ and $D(C)$ for all $D \in \mathcal{D}$. It follows by Fact 1.24 that $\operatorname{bdn}(A / D(A) \times D(C))=\operatorname{bdn}(A / D(A))+\operatorname{bdn}(D(C))$. The interpretable (and surjective) map

$$
\begin{array}{ccc}
\rho_{D} \times \nu: \quad \nu^{-1}(D(C)) & \rightarrow & A / D(A) \times D(C) \\
b & \mapsto & \left(\rho_{D}(b), \nu(b)\right)
\end{array}
$$

gives us that

$$
\operatorname{bdn}(B) \geq \max _{D \in \mathcal{D}}(\operatorname{bdn}(A / D(A))+\operatorname{bdn}(D(C)))
$$

It remains to show that $\operatorname{bdn} B \leq \max _{D \in \mathcal{D}}(\operatorname{bdn}(A / D(A))+\operatorname{bdn}(D(C)))$. By Remark 2.3, we may assume that $\operatorname{bdn}(A)$ and $\operatorname{bdn}(C)$ are both finite. As $A$ is infinite, $\operatorname{bdn} A \geq 1$. If $\operatorname{bdn}(\mathcal{M})=1$, the equality is clear. Assume that $\operatorname{bdn}(\mathcal{M})>1$ and let $P(x)=\left\{\phi_{i}\left(x, y_{i}\right),\left(a_{i, j}\right)_{j<\omega}, k_{i}\right\}_{i<M}$ be an inppattern of finite depth $M \geq 2$, with $\left(a_{i, j}\right)_{i, j}$ mutually indiscernible and $|x|=1$. We need to show that $M \leq \operatorname{bdn}(A / D(A))+\operatorname{bdn}(D(C))$ for some $D \in \mathcal{D}$. If $x$ is a variable in the sort $A$ (resp. in the sort $C), P(x)$ is an inp-pattern in $A$ (resp. $C$ ) of depth bounded by $\operatorname{bdn}(A)$ (resp. $\operatorname{bdn}(C)$ ) by purity (Corollary 1.75). Then, the inequality holds if we take $D \equiv x=x$ (respect $D \equiv x=0$ ).

Assume $x$ is a variable in the sort $B$. Consider a line $\left\{\phi(x, y),\left(a_{j}\right)_{j<\omega}\right\}$ of $P(x)$ (we drop the index $i<M$ for the sake of clarity). By Fact 1.74, and by the fact that one can "eliminate" disjunctions in inp-patterns (see Lemma 1.23), we may assume that the formula $\phi\left(x, a_{j}\right)$ is of the form

$$
\begin{align*}
& \phi_{A}\left(\rho_{n_{0}}\left(t^{0}\left(x, \beta_{j}\right)\right), \ldots, \rho_{n_{s-1}}\left(t^{s-1}\left(x, \beta_{j}\right)\right), \alpha_{j}\right)  \tag{5}\\
& \quad \wedge \phi_{C}\left(\nu\left(r^{0}\left(x, \beta_{j}\right)\right), \ldots, \nu\left(r^{k-1}\left(x, \beta_{j}\right)\right), \gamma_{j}\right), \tag{6}
\end{align*}
$$

where $\phi_{A}$ is an $A$-formula, $\phi_{C}$ a $C$-formula, and for $j<\omega, \alpha_{j} \in A, \beta_{j} \in B, \gamma_{j} \in C$ are parameters, $s, k, n_{0}, n_{1}, \ldots, n_{s-1} \in \mathbb{N}$ and the $t^{l}$ s and $r^{l}$ 's are terms in the group language (one needs to keep in mind that $s, k, n_{l}, t^{l}, r^{l}, \beta_{j}, \alpha_{j}$ and $\gamma_{j}$ depend on the line $\left.i\right)$. Also, notice that $\rho_{n_{0}}\left(t^{0}\left(x, \beta_{j}\right)\right) \neq 0 \in A / n_{0} A$ implies $\nu\left(t^{0}\left(x, \beta_{j}\right)\right) \in n_{0} C$ (a formula of the form (6)). We may write

$$
\begin{aligned}
\phi_{A}\left(\rho _ { n _ { 0 } } \left(t^{0}\right.\right. & \left.\left.\left(x, \beta_{j}\right)\right), \ldots, \rho_{n_{s-1}}\left(t^{s-1}\left(x, \beta_{j}\right)\right), \alpha_{j}\right) \simeq \\
& \left(\phi_{A}\left(\rho_{n_{0}}\left(t^{0}\left(x, \beta_{j}\right)\right), \ldots, \rho_{n_{s-1}}\left(t^{s-1}\left(x, \beta_{j}\right)\right), \alpha_{j}\right) \wedge \nu\left(t^{0}\left(x, \beta_{j}\right)\right) \in n_{0} C\right) \\
& \bigvee\left(\phi_{A}\left(0, \rho_{n_{1}}\left(t^{1}\left(x, \beta_{j}\right)\right), \ldots, \rho_{n_{s-1}}\left(t^{s-1}\left(x, \beta_{j}\right)\right), \alpha_{j}\right) \wedge \nu\left(t^{0}\left(x, \beta_{j}\right)\right) \notin n_{0} C\right)
\end{aligned}
$$

and use once again Lemma 1.23 to eliminate the disjunction. We have then one of these two cases:

- either the first part of the disjunction remains, and we have an occurrence of the term $\rho_{n_{0}}\left(t^{0}\left(x, \beta_{j}\right)\right)$ and of the atomic formula $\nu\left(t^{0}\left(x, \beta_{j}\right)\right) \in n_{0} C$,
- or the second part of the disjunction remains, and the term $\rho_{n_{0}}\left(t^{0}\left(x, \beta_{j}\right)\right)$ does not occur anymore, as it has been replaced by 0 . In that case, we reset the notation: $\rho_{n_{1}}\left(t^{1}\left(x, \beta_{j}\right)\right)$ becomes our new first term ' $\rho_{n_{0}}\left(t^{0}\left(x, \beta_{j}\right)\right)$ '; and $\rho_{n_{2}}\left(t^{2}\left(x, \beta_{j}\right)\right)$ becomes ' $\rho_{n_{1}}\left(t^{1}\left(x, \beta_{j}\right)\right.$ )' etc.

It follows that we can assume that $\phi\left(x, a_{j}\right)$ (or more specifically, the formula $\phi_{C}\left(\nu\left(t^{0}\left(x, \beta_{j}\right), \ldots, \nu\left(t^{s-1}\left(x, \beta_{j}\right)\right), \gamma_{j}\right)\right)$ implies $t^{0}\left(x, \beta_{j}\right) \in \nu^{-1}\left(n_{0} C\right)$. We do the same for all terms $t^{l}\left(x, \beta_{j}\right), l<s$. This means in particular that the list of terms $\left\{t^{l}\right\}_{l<s}$ is included in $\left\{r^{l}\right\}_{l<k}$.

Let $M^{\prime} \geq 0$ be the number of rows such that

$$
\left\{\phi_{C}\left(\nu\left(r^{0}\left(x, \beta_{j}\right)\right), \ldots, \nu\left(r^{k-1}\left(x, \beta_{j}\right)\right), \gamma_{j}\right)\right\}_{j<\omega}
$$

is consistent. Without loss, they are the $M^{\prime}$ first rows of the pattern $P(x)$, and we denote by $P^{\prime}(x)$ the sub-pattern consisting of these rows. For now, we work with the sub-pattern $P^{\prime}(x)$.

Terms $t\left(x, \beta_{j}\right)$ in the group language are of the form $k x+m \cdot \beta_{j}$, with $k \in \mathbb{N}, m \in \mathbb{N}^{\left|\beta_{j}\right|}$.
Claim 2. Assume that, in a line $\left\{\phi(x, y),\left(a_{j}\right)_{j<\omega}\right\}$ of $P^{\prime}(x)$, a term $\rho_{n}\left(k x+m \cdot \beta_{j}\right)$ occurs. Then $\nu\left(m \cdot \beta_{j}\right) \bmod n C$ is constant for all $j<\omega$.

Proof. Assume not. By indiscernibility, $\nu\left(m \cdot \beta_{j}\right)$ are in distinct classes modulo $n C$. As

$$
\phi_{C}\left(\nu\left(r^{0}\left(x, \beta_{j}\right)\right), \ldots, \nu\left(r^{s-1}\left(x, \beta_{j}\right)\right), \gamma_{j}\right) \vdash \nu\left(k x+m \cdot \beta_{j}\right) \in n C,
$$

the $\phi_{C}$ part of the line

$$
\left\{\phi_{C}\left(\nu\left(r^{0}\left(x, \beta_{j}\right)\right), \ldots, \nu\left(r^{s-1}\left(x, \beta_{j}\right)\right), \gamma_{j}\right)\right\}_{j<\omega}
$$

is 2 -inconsistent, contradicting the fact we chose one of the first $M^{\prime}$ lines.
Claim 3 (Main claim). We may assume that all formulas in $P^{\prime}(x)$ are of the form

$$
\begin{aligned}
& \phi_{A}\left(\rho_{n_{0}}\left(k^{0}(x-d)\right), \ldots, \rho_{n_{s-1}}\left(k^{s-1}(x-d)\right), \alpha_{j}\right) \\
& \wedge \phi_{C}\left(\nu(x-d), \gamma_{j}\right) .
\end{aligned}
$$

for certain integer $n_{0}, \cdots, n_{s-1}, k^{0}, \cdots, k^{s-1}$ and a certain parameter $d \in B$.
Proof. Take any realisation $d$ of the first column:

$$
d \models\left\{\phi_{i}\left(x, a_{i, 0}\right)\right\}_{i<M} .
$$

We fix an $i<M^{\prime}$, and consider the $i^{\text {th }} \operatorname{line}\left\{\phi\left(x, a_{j}\right)\right\}_{j<\omega}$, (again, we drop the index $i$ for a simpler notation).

Step 1: We may assume that all terms $t^{l}\left(x, \beta_{j}\right)=k^{l} x+m^{l} \cdot \beta_{j}, l<s$ are of the form $k^{l}(x-d)$. We change all terms one by one, starting with $t^{0}\left(x, \beta_{j}\right)=k^{0} x+m^{0} \cdot \beta_{j}$. We write

$$
\rho_{n_{0}}\left(k^{0} x+m^{0} \cdot \beta_{j}\right)=\rho_{n_{0}}\left(k^{0}(x-d)+k^{0} d+m^{0} \cdot \beta_{j}\right) .
$$

Replace it by $\rho_{n_{0}}\left(k^{0}(x-d)\right)+\rho_{n_{0}}\left(k^{0} d+m^{0} \cdot \beta_{j}\right)$. This doesn't change the formula $\phi\left(x, a_{j}\right)$ as $\phi\left(x, a_{j}\right) \vdash$ $\nu\left(k^{0}(x-d)\right) \in n_{0} C$. Indeed, $\phi\left(x, a_{j}\right) \vdash \nu\left(k^{0} x+m^{0} \cdot \beta_{j}\right) \in n_{0} C$ and $\nu\left(k^{0} d+m^{0} \cdot \beta_{j}\right) \in n_{0} C$ since $d$ is a solution of the first column and $\nu\left(m^{0} \cdot \beta_{j}\right) \bmod n_{0} C$ is constant by the previous claim. Then, $\rho_{n_{0}}\left(k^{0} d-m^{0} \cdot \beta_{j}\right)$ is seen as a parameter in $A / n_{0} A$, and it is added to $\alpha_{j}$. We do the same for all terms $t^{l}\left(x, \beta_{j}\right), l<s$.

Step 2: We may assume that all terms $r^{l}\left(x, \beta_{j}\right)=k^{l} x+m^{l} \cdot \beta_{j}, l<k$, are of the form $x-d$. This is immediate, as $\nu$ is a morphism. Indeed, replace $\nu\left(r^{l}\left(x, \beta_{j}\right)\right)$ by $k^{l} \nu(x-d)+\nu\left(k^{l} d+m^{l} \cdot \beta_{j}\right)$, where $\nu\left(k^{l} d+m^{l} \cdot \beta_{j}\right)$ is seen as a parameter in $C$, and is added to the parameters $\gamma_{j}$. In other words, we may assume that the formula in the $i$ th line of the pattern $P^{\prime}(x)$ has the required form.

Let us recall finally that $d$ has been chosen independently of the row. We can apply these steps for all rows $i<M^{\prime}$.

Let $D$ be the conjunction of the formulas $\exists y k x=n y$ for each term $\rho_{n}(k(x-d))$ which occurs in any line of $P^{\prime}(x)$. Notice that $\left\{\phi_{i}\left(x, a_{i, f(i)}\right)\right\}_{i<M^{\prime}}$ implies that $\nu(x-d) \in D(C)$ for any choice of function $f: M^{\prime} \rightarrow \omega$. If $x_{A}$ is a variable in $A$, notice also that the truth value of the formula

$$
\phi_{A}\left(\pi_{n_{0}}\left(k^{0} x_{A}\right), \ldots, \pi_{n_{s-1}}\left(k^{s-1} x_{A}\right), \alpha_{j}\right)
$$

evaluated in $a \in A$ is independent of the class of $a$ modulo $D(A)$. In other words, it interprets a definable set in $A / D(A)$.

As now the pattern $P^{\prime}(x)$ is centralised, one can remark that for every line

$$
\left\{\phi\left(x, a_{j}\right) \equiv \phi_{A}\left(\rho_{n_{0}}\left(k^{0}(x-d)\right), \ldots, \rho_{n_{s-1}}\left(k^{s-1}(x-d)\right), \alpha_{j}\right) \wedge \phi_{C}\left(\nu(x-d), \gamma_{j}\right)\right\}_{j<\omega}
$$

at most one of the following sets:

$$
\begin{equation*}
\left\{\phi_{A}\left(\pi_{n_{0}}\left(k^{0} x_{A}\right), \ldots, \pi_{n_{s-1}}\left(k^{s-1} x_{A}\right), y_{A}\right),\left(\alpha_{j}\right)_{j<\omega}\right\} \tag{A}
\end{equation*}
$$

where $\left|y_{A}\right|=\left|\alpha_{j}\right|$ and $x_{A}$ is a variable in $A$ or

$$
\begin{equation*}
\left\{\phi_{C}\left(x_{C}, y_{C}\right),\left(\gamma_{j}\right)_{j<\omega}\right\} \tag{C}
\end{equation*}
$$

(where $\left|y_{C}\right|=\left|\gamma_{j}\right|$ ) is consistent. Indeed, this follows immediately from the fact that in the monster model, the sequence splits and $B \simeq A \times C$. By definition of $P^{\prime}(x)$, we deduce that $\left(L_{A}\right)$ - the $\phi_{A}$-part of the line- is inconsistent. Now, take a path $f: M^{\prime} \rightarrow \omega$ and $b$ a solution of

$$
\left\{\phi_{i}\left(x, a_{i, f(i)}\right)\right\}_{i<M^{\prime}}
$$

As $b-d \in D(C)$, there is an $a$ such that $\iota(a)+D(B)=b-d+D(b)$. This give us that

$$
a \models\left\{\phi_{i, A}\left(\pi_{n_{0}}\left(k^{0} x_{A}\right), \ldots, \pi_{n_{s-1}}\left(k^{s-1} x_{A}\right), \alpha_{i, f(i)}\right)\right\}_{i<M^{\prime}}
$$

This show that every path of the following pattern $P_{A}^{\prime}\left(x_{A}\right)$ is consistent:

$$
P_{A}^{\prime}\left(x_{A}\right):=\left\{\phi_{i, A}\left(\pi_{n_{0}}\left(k^{0} x_{A}\right), \ldots, \pi_{n_{s-1}}\left(k^{s-1} x_{A}\right), y_{A}\right),\left(\alpha_{i, j}\right)_{j<\omega}\right\}_{i<M^{\prime}}
$$

Thus, this is an inp-pattern in $A$ and more precisely, it interprets an inp-pattern in $A / D(A)$ by the remark above. This means that $M^{\prime} \leq \operatorname{bdn}(A / D(A))$.

Then, the $\phi_{C}$-part $\left(L_{C}\right)$ of any of the $M-M^{\prime}$ last lines of $P(x)$ is inconsistent by definition of $P^{\prime}(x)$ and $M^{\prime}$. One gets as well an inp-pattern of depth $M-M^{\prime}$ in $C$. As any realisation $r$ of $P^{\prime}(x)$ (in particular, of $P(x))$ satisfies $\nu(r-d) \in D(C)$, one gets actually an inp-pattern in $D(C)$. It follows that $M-M^{\prime} \leq \operatorname{bdn}(D(C))$. At the end, we get that $M \leq \operatorname{bdn}(A / D(A))+\operatorname{bdn}(D(C))$.

To conclude, we treat the particular cases. Assume that $C$ has finite $n$-torsion for every $n \in \mathbb{N}^{\star}$. Then infinite subgroups in $\mathcal{D}(C):=\{D(C) \mid D \in \mathcal{D}\}$ are of the form $n C$. We deduce that

$$
\operatorname{bdn} \mathcal{M}=\max _{n \in \mathbb{N}}(\operatorname{bdn}(A / n A)+\operatorname{bdn}(n C))
$$

Similarly assume that $A / n A$ is finite for all $n \in \mathbb{N}^{\star}$. Then $\operatorname{bdn}\left(A / n A_{[k]}\right)$ is equal to $\operatorname{bdn}\left(A / A_{[k]}\right)$. Indeed, this can be deduced from Fact 1.26 , as we have the exact sequence

$$
0 \rightarrow A_{[k]} / n A_{[k]} \rightarrow A / n A_{[k]} \rightarrow A / A_{[k]} \rightarrow 0
$$

and as $A_{[k]} / n A_{[k]}$ is finite. Since $\operatorname{bdn}\left(C_{[k]}\right) \geq \operatorname{bdn}\left(n C_{[k]}\right)$ and $A / A_{[k]} \simeq k A$, we may deduce that

$$
\operatorname{bdn} \mathcal{M}=\max _{k \in \mathbb{N}}\left(\operatorname{bdn}(k A), \operatorname{bdn}\left(C_{[k]}\right)\right)
$$

### 2.2 Application

As main application of Theorem 2.2, we will deduce Theorem 3.12. This is the aim of the next section. For now, we want to emphasise the advantage of working resplendently by giving one straightforward generalisation of Theorem 2.2.
Corollary 2.5. Let $\mathcal{M}$ be an exact sequence of ordered abelian groups

$$
0 \longrightarrow A \xrightarrow{\iota} B \xrightarrow{\nu} C \longrightarrow 0,
$$

where $(A,<)$ is a convex subgroup of $(C,<)$. We consider it as a three sorted structure, with a structure of ordered abelian group for each sort, and function symbols for $\iota$ and $\nu$. Then, we have:

$$
\begin{aligned}
\operatorname{bdn} \mathcal{M} & =\max _{n \in \mathbb{N}}(\operatorname{bdn}(A / n A)+\operatorname{bdn}(n C)) \\
& =\max \left(\operatorname{bdn}(A), \max _{n \in \mathbb{N}_{\star}}(\operatorname{bdn}(A / n A))+\operatorname{bdn}(C)\right) .
\end{aligned}
$$

Proof. As $C$ is torsion free, $\iota(A)$ is pure in $B$. As for $b \in B, b>0$ if and only if $\nu(b)>0$ or $\nu(b)=0$ and $\iota^{-1}(b)>0, \mathcal{M}$ is an $\{A\}-\{C\}$-enrichment of a short exact sequence of abelian groups. It remains to apply Theorem 2.2. Notice that for all $n>0$, we have $\operatorname{bdn}(n C)=\operatorname{bdn}(C)$ (as the multiplication by $n$ in $C$ is a definable injective morphism).

## 3 Burden of Henselian valued fields

We compute the burden of benign Henselian valued fields and of unramified mixed characteristic Henselian valued fields with perfect residue field in terms of the burden of the value group and that of the residue field. The first subsection is common for both cases and treats the reduction from the valued field to the sort RV (resp. the sorts $R V_{<\omega}$ ). For the reduction to the value group and residue field, we treat (separately) the case of benign Henselian valued fields in Subsection 3.2 and the case of unramified mixed characteristic henselian valued fields with perfect residue field in Subsection 3.3. They are both deduced from the computation of burden in short exact sequences of abelian groups (Section 2).

### 3.1 Reductions to RV and $\mathrm{RV}_{<\omega}$

We compute here the burden of Henselian valued fields of characteristc 0 in terms of burden of $R V_{<\omega}$. As we explained in Paragraph 1.2.1, mixed characteristic Henselian valued fields satisfy (EQ) $\mathrm{RV}_{<\omega}$ (quantifier elimination relative to the union of sorts $\mathrm{RV}_{<\omega}$.), but do not satisfy in general (EQ) RV (elimination of quantifiers relative to RV). Our result includes naturally the case of unramified mixed Henselian valued fields, and also equicharacteristic 0 Henselian valued fields. In the former case, we have a computation of the burden in term of the burden of RV, as in equicharacteristic 0 the structures $R V$ and $R V_{<\omega}$ can be identified (Remark 1.46). In fact, the proof that we are going to present can be adapted for all benign valued fields.

### 3.1.1 Reductions

The aim of this paragraph is to prove the following:
Theorem 3.1. Let $K$ be a Henselian valued field of characteristic $(0, p), p \geq 0$. Let $M$ be a positive integer and assume $\mathcal{K}$ is of burden $M$. Then, the sort $\mathrm{RV}_{<\omega}$ with the induced structure is also of burden M. In particular, $\mathcal{K}$ is inp-minimal if and only if $\mathrm{RV}_{<\omega}$ is inp-minimal.

The demonstration (below) follows Chernikov and Simon's proof for the case of equicharacteristic 0 and burden 1 (see [11]). As we said, this statement also cover the case of equicharacteristic 0 . One can also generalise the proof for infinite burden (see Corollary 3.3 for details). A careful reading of the proof shows that one only uses properties $(E Q)_{R V_{<\omega}}$ and (Lin) $\mathrm{RV}_{<\omega}$. Of course, the proof can be written for equicaracteristic 0 fields only (it becomes simpler), and then it only uses Property (EQ) RV and (Lin) RV. As algebraically maximal Kaplansky valued fields and algebraically closed valued fields satisfy these property, we obtain in fact:

Theorem 3.2. Let $\mathcal{K}$ be a benign Henselian valued field. Let $M$ be a positive integer and assume $\mathcal{K}$ is of burden $M$. Then, the sort RV with the induced structure is also of burden $M$. In particular, $\mathcal{K}$ is inp-minimal if and only if RV is inp-minimal.

Proof of Theorem 3.1. We denote by $\overline{\mathbb{Z}}$ the set of natural numbers with extremal points $\mathbb{Z} \cup\{ \pm \infty\}$. Let $\left\{\tilde{\phi}_{i}\left(x, y_{i}\right),\left(c_{i, j}\right)_{j \in \overline{\mathbb{Z}}}, k_{i}\right\}_{i<M}$ be an inp-pattern in $\mathcal{K}$ of finite depth $M \geq 2$ with $|x|=1$, where $c_{i, j}=a_{i, j} \mathbf{b}_{i, j} \in K^{k_{1}} \times \mathrm{RV}_{<\omega}^{k_{2}}$. Notice that the set of indices is $\overline{\mathbb{Z}}$, as we will make use of one of the extreme elements $\left.\left\{a_{i,-\infty}, a_{i,+\infty}\right\}\right)$ later. We have to find an inp-pattern of depth $M$ in $\mathrm{RV}_{<\omega}$. Without loss of generality, we take $\left(c_{i, j}\right)_{i, j}$ mutually indiscernible. By Fact 1.64 and mutual indiscernibility, we can assume the formulas $\tilde{\phi}_{i}$ are of the form

$$
\tilde{\phi}_{i}\left(x, c_{i, j}\right)=\phi_{i}\left(\operatorname{rv}_{\delta_{n}}\left(x-a_{i, j ; 1}\right), \ldots, \operatorname{rv}_{\delta_{n}}\left(x-a_{i, j ; k_{1}}\right) ; \mathbf{b}_{i, j}\right),
$$

for some integer $n$ and where $\phi_{i}$ are $\mathrm{RV}_{<\omega}$-formulas. Also recall that $\delta_{n}$ denotes the value $\operatorname{val}\left(p^{n}\right)$. The arguments inside symbols $\mathrm{rv}_{\delta_{n}}$ are linear terms in $x$. In some sense, difficulties coming from the field structure have been already treated and it only remains to deal with the structure coming from the valuation.

Let $d \equiv\left\{\phi_{i}\left(\operatorname{rv}_{\delta_{n}}\left(x-a_{i, 0 ; 1}\right), \ldots, \operatorname{rv}_{\delta_{n}}\left(x-a_{i, 0 ; k_{1}}\right) ; \mathbf{b}_{i, 0}\right)\right\}_{i<M}$ be a solution of the first column. Before we give a general idea of the proof, let us reduce to the case where only one term $\operatorname{rv}_{\delta_{n}}\left(x-a_{i, j}\right)$ occurs in the formula $\tilde{\phi}_{i}$.
Claim 4. We may assume that for all $i<M, \tilde{\phi}_{i}\left(x, c_{i, j}\right)$ is of the form $\phi_{i}\left(\operatorname{rv}_{\delta_{n}}\left(x-a_{i, j ; 1}\right) ; \mathbf{b}_{i, j}\right)$, i.e. $\left|a_{i, j}\right|=k_{1}=1$.

Proof. We will first replace the formula $\tilde{\phi}_{0}\left(x, c_{0, j}\right)$ by a new one with an extra parameter.
By Lemma 1.50, at least one of the following two cases occurs

1. $\mathrm{WD}_{\delta_{n}}\left(\operatorname{rv}_{\delta_{n}}\left(d-a_{0,0 ; 1}\right), \operatorname{rv}_{\delta_{n}}\left(a_{0,0 ; 1}-a_{0,0 ; 2}\right)\right)$ or
2. $\mathrm{WD}_{\delta_{n}}\left(\operatorname{rv}_{\delta_{n}}\left(d-a_{0,0 ; 2}\right), \operatorname{rv}_{\delta_{n}}\left(a_{0,0 ; 2}-a_{0,0 ; 1}\right)\right)$.

According to the case, we respectively define a new formula $\psi_{0}\left(x, c_{0, j}{ }^{\wedge} \operatorname{rv}_{\delta_{n}}\left(a_{0, j ; 2}-a_{0, j ; 1}\right)\right)$ by:
1.

$$
\begin{aligned}
& \phi_{0}\left(\operatorname{rv}_{\delta_{n}}\left(x-a_{0, j ; 1}\right), \operatorname{rv}_{\delta_{n}}\left(x-a_{0, j ; 1}\right)+\operatorname{rv}_{\delta_{n}}\left(a_{0, j ; 1}-a_{0, j ; 2}\right), \operatorname{rv}_{\delta_{n}}\left(x-a_{0, j ; 3}\right), \ldots\right. \\
&\left.\operatorname{rv}_{\delta_{n}}\left(x-a_{0, j ; k}\right) ; \mathbf{b}_{0, j}\right) \wedge \mathrm{WD}_{\delta_{n}}\left(\operatorname{rv}_{\delta_{n}}\left(x-a_{0, j ; 1}\right), \operatorname{rv}_{\delta_{n}}\left(a_{0, j ; 1}-a_{0, j ; 2}\right)\right),
\end{aligned}
$$

2. 

$$
\begin{aligned}
\phi_{0}\left(\operatorname{rv}_{\delta_{n}}\left(x-a_{0, j ; 2}\right)+\operatorname{rv}_{\delta_{n}}\left(a_{0, j ; 2}-a_{0, j ; 1}\right), \operatorname{rv}_{\delta_{n}}( \right. & \left.\left.x-a_{0, j ; 2}\right), \ldots, \operatorname{rv}_{\delta_{n}}\left(x-a_{0, j ; k}\right) ; \mathbf{b}_{0, j}\right) \\
& \wedge \mathrm{WD}_{\delta_{n}}\left(\operatorname{rv}_{\delta_{n}}\left(x-a_{0, j ; 2}\right), \operatorname{rv}_{\delta_{n}}\left(a_{0, j ; 2}-a_{0, j ; 1}\right)\right)
\end{aligned}
$$

We will prove that the pattern where $\tilde{\phi}_{0}$ is replaced by $\psi_{0}$ :

$$
\left\{\psi_{0}\left(x, y_{0} \wedge z\right),\left(c_{0, j}^{\wedge} \operatorname{rv}_{\delta_{n}}\left(a_{0, j ; 2}-a_{0, j ; 1}\right)\right)_{j \in \overline{\mathbb{Z}}}, k_{0}\right\} \cup\left\{\tilde{\phi}_{i}\left(x, y_{i}\right),\left(c_{i, j}\right)_{j \in \overline{\mathbb{Z}}}, k_{i}\right\}_{1 \leq i<M}
$$

is also an inp-pattern. First note that we have added $\operatorname{rv}_{\delta_{n}}\left(a_{0, j ; 2}-a_{0, j ; 1}\right)$ to the parameters $\mathbf{b}_{0, j}$, and it still forms a mutually indiscernible array. Clearly, $d$ is still a realisation of the first column:

$$
d \equiv\left\{\psi_{0}\left(x, c_{0,0} \wedge \operatorname{rv}_{\delta_{n}}\left(a_{0,0 ; 2}-a_{0,0 ; 1}\right)\right)\right\} \cup\left\{\tilde{\phi}_{i}\left(x, c_{i, 0}\right) \mid 1 \leq i<M\right\} .
$$

By mutual indiscernibility of the parameters, every path is consistent. Since $\psi_{0}(\mathcal{K}) \subseteq \tilde{\phi}_{0}(\mathcal{K})$, inconsistency of the first row is also clear. By induction, it is clear that we may assume that $\phi_{0}$ is of the desired form. We can do the same for all formulas $\phi_{i}, 0<i<M$.

If the array $\left(a_{i, j}\right)_{i<M, j<\omega}$ is constant equal to some $a \in K$, then we obviously get an inp-pattern of depth $M$ in $\mathrm{RV}_{<\omega}:\left\{\phi_{i}\left(\mathbf{x}, z_{i}\right),\left(\mathbf{b}_{i, j}\right)_{j \in \overline{\mathbb{Z}}}, k_{i}\right\}_{i<M}$, where $\mathbf{x}$ is a variable in $\mathrm{RV}_{\delta_{n}}$ (such a pattern is said to be centralised ). Indeed, consistency of the path is clear. If a row is satisfied by some $\mathbf{d} \in \operatorname{RV}_{\delta_{n}}$ , any $d \in K$ such that $\operatorname{rv}_{\delta_{n}}(d-a)=\mathbf{d}$ will satisfy the corresponding row of the initial inp-pattern, which is absurd. Hence, the rows are inconsistent.

The idea of the proof is to reduce the general case (where the $a_{i, j}$ 's are distinct) to this trivial case by the same method as above: removing the parameters $a_{i, j} \in K$, adding new parameters from $\mathrm{RV}_{<\omega}$ to $\mathbf{b}_{i, j}$ and adding a term of the form $\mathrm{WD}\left(\operatorname{rv}(x-a), \operatorname{rv}\left(a-a_{i, j}\right)\right)$. The main challenge is to find a suitable $a \in K$ for a center.

Recall that $d \models\left\{\phi_{i}\left(\operatorname{rv}_{\delta_{n}}\left(x-a_{i, 0}\right) ; b_{i, 0}\right)\right\}_{i<M}$ is a solution to the first column.
Claim 5. For all $j<\omega$, and $i, k<M$ with $k \neq i$, we have $\operatorname{val}\left(d-a_{i, j}\right) \leq \operatorname{val}\left(d-a_{k, 0}\right)+\delta_{n}$.
Proof. Assume not: for some $j<\omega$, and $i, k<M$ with $k \neq i$ :

$$
\operatorname{val}\left(d-a_{i, j}\right)>\operatorname{val}\left(d-a_{k, 0}\right)+\delta_{n}
$$

Then, $\operatorname{rv}_{\delta_{n}}\left(a_{i, j}-a_{k, 0}\right)=\operatorname{rv}_{\delta_{n}}\left(d-a_{k, 0}\right)$. By mutual indiscernibility, we have

$$
a_{i, j} \models\left\{\phi_{k}\left(\operatorname{rv}_{\delta_{n}}\left(x-a_{k, l}\right) ; \mathbf{b}_{k, l}\right)\right\}_{l<\omega} .
$$

This contradicts inconsistency of the row $k$.
In particular, for all $i, k<M$, we have $\left|\operatorname{val}\left(d-a_{k, 0}\right)-\operatorname{val}\left(d-a_{i, 0}\right)\right| \leq \delta_{n}$. For $i<M$, let us denote $\gamma_{i}:=\operatorname{val}\left(d-a_{i, 0}\right)$ and let $\gamma$ be the minimum of the $\gamma_{i}$ 's. By definition, we have the following for all $i, k<M$ :

$$
\operatorname{val}\left(a_{i, 0}-a_{k, 0}\right) \geq \min \left\{\operatorname{val}\left(d-a_{i, 0}\right), \operatorname{val}\left(d-a_{k, 0}\right)\right\} \geq \gamma
$$

The following claim give us a correct centre $a$.
Claim 6. We may assume that there is $i<M$ such that for all $k<M$, the following holds:

$$
\gamma_{k}=\operatorname{val}\left(d-a_{k, 0}\right) \leq \min \left\{\operatorname{val}\left(d-a_{i, \infty}\right), \operatorname{val}\left(a_{i, \infty}-a_{k, 0}\right)\right\}+\delta_{n}
$$

In particular, by Proposition 1.48, we have:

$$
\mathrm{WD}_{\delta_{n}}\left(\mathrm{rv}_{2 \delta_{n}}\left(d-a_{i, \infty}\right), \mathrm{rv}_{2 \delta_{n}}\left(a_{i, \infty}-a_{k, 0}\right)\right)
$$

Proof. By Remark 1.49, it is enough to find $i<M$ such that the following holds for all $k<M$ :

$$
\gamma_{k} \leq \operatorname{val}\left(d-a_{i, \infty}\right)+\delta_{n} \quad \text { or } \quad \gamma_{k} \leq \operatorname{val}\left(a_{i, \infty}-a_{k, 0}\right)+\delta_{n}
$$

We will actually find $i$ such that one of the following holds:

1. $\gamma_{k} \leq \operatorname{val}\left(d-a_{i, \infty}\right)+\delta_{n}$ for all $k<M$
2. $\gamma_{k} \leq \operatorname{val}\left(a_{i, \infty}-a_{k, 0}\right)+\delta_{n}$ for all $k<M$

The first case will correspond to Case A, the second to Case B.
Case A: There are $0 \leq i, k<M$ with $i \neq k$ such that $\operatorname{val}\left(a_{i, j}-a_{k, l}\right)$ is constant for all $j, l \in \omega$, equal to some $\epsilon$. Note that ( $\star$ ) gives $\epsilon \geq \gamma$.


Then, we have:

$$
\operatorname{val}\left(d-a_{i, \infty}\right) \geq \min \left\{\operatorname{val}\left(d-a_{i, 0}\right), \operatorname{val}\left(a_{i, 0}-a_{i, \infty}\right)\right\} \geq \gamma
$$

Indeed, $\operatorname{val}\left(a_{i, 0}-a_{i, \infty}\right) \geq \min \left\{\operatorname{val}\left(a_{i, 0}-a_{k, 0}\right), \operatorname{val}\left(a_{k, 0}-a_{i, \infty}\right)\right\}=\epsilon \geq \gamma$. Hence, we have for every $0 \leq l<M$ :

$$
\operatorname{val}\left(d-a_{i, \infty}\right)+\delta_{n} \geq \gamma+\delta_{n} \geq \operatorname{val}\left(d-a_{l, 0}\right)=\gamma_{l}
$$

Case B: For all $0 \leq i, k<M$ with $i \neq k,\left(\operatorname{val}\left(a_{i, j}-a_{k, l}\right)\right)_{j, l}$ is not constant.
By Lemma $1.40(2)$ and Lemma 1.41, there is $i<M$ such that for every $k<M$ and $k \neq i,\left(a_{k, l}\right)_{l<\omega} \Rightarrow a_{i, 0}$ or $\left(a_{k,-l}\right)_{l<\omega} \Rightarrow a_{i, 0}$. If needed, one can flip the indices and assume that for all $k \neq i,\left(a_{k, l}\right)_{l<\omega} \Rightarrow a_{i, 0}$. Note that only $\left(a_{i, j}\right)_{j}$ could be a fan in this case.


Then we have

$$
\operatorname{val}\left(a_{k, 0}-a_{i, \infty}\right)=\operatorname{val}\left(a_{k, 0}-a_{i, 0}\right) \geq \gamma
$$

since $\left(a_{k, l}\right) \Rightarrow a_{i, \infty}$ as well. So

$$
\operatorname{val}\left(a_{k, 0}-a_{i, \infty}\right)+\delta_{n} \geq \gamma_{k}
$$

It remains to prove the inequality for $k=i$. Take $l \neq i, l<M$. We have:

$$
\operatorname{val}\left(a_{i, 0}-a_{i, \infty}\right) \geq \min \left\{\operatorname{val}\left(a_{i, 0}-a_{l, 0}\right), \operatorname{val}\left(a_{l, 0}-a_{i, \infty}\right)\right\} \geq \gamma
$$

Hence, $\operatorname{val}\left(a_{i, 0}-a_{i, \infty}\right)+\delta_{n} \geq \gamma_{i}$.

Assume $i=0$ satisfies the conclusion of the previous claim. For every $k<M$, we have

$$
\mathrm{WD}_{\delta_{n}}\left(\operatorname{rv}_{2 \delta_{n}}\left(d-a_{0, \infty}\right), \operatorname{rv}_{2 \delta_{n}}\left(a_{0, \infty}-a_{k, 0}\right)\right)
$$

Set $\tilde{b}_{i, j}:=b_{i, j}{ }^{\wedge} \operatorname{rv}_{2 \delta_{n}}\left(a_{0, \infty}-a_{i, j}\right)$ for $i<M, j<\omega$ and

$$
\psi_{i}\left(\tilde{x}, \tilde{b}_{i, j}\right):=\phi_{i}\left(\operatorname{rv}_{\delta_{n}}\left(\tilde{x}+\operatorname{rv}_{2 \delta_{n}}\left(a_{0, \infty}-a_{i, j}\right)\right) ; b_{i, j}\right) \wedge \mathrm{WD}_{\delta_{n}}\left(\tilde{x}, \operatorname{rv}_{2 \delta_{n}}\left(a_{0, \infty}-a_{i, j}\right)\right)
$$

where $\tilde{x}$ is a variable in $\mathrm{RV}_{2 \delta_{n}}$.
This is an inp-pattern. Indeed, clearly, $\operatorname{rv}_{2 \delta_{n}}\left(d-a_{0, \infty}\right) \models\left\{\psi_{i}\left(\tilde{x}, \tilde{b}_{i, 0}\right)\right\}_{i<M}$. By mutual indiscernibility of $\left(\tilde{b}_{i, j}\right)_{i<M, j<\omega}$, every path is consistent. It remains to show that, for every $i<M,\left\{\psi_{i}\left(\tilde{x}, \tilde{b}_{i, j}\right)\right\}_{j<\omega}$ is inconsistent. Assume there is $\alpha^{\star} \models\left\{\psi_{i}\left(\tilde{x}, \tilde{b}_{i, j}\right)\right\}_{j<\omega}$ for some $i<M$, and let $d^{\star}$ be such that $\operatorname{rv}_{2 \delta_{n}}\left(d^{\star}-a_{0, \infty}\right)=\alpha^{\star}$. Then, since $\mathrm{WD}_{\delta_{n}}\left(\alpha^{\star}, \operatorname{rv}_{2 \delta_{n}}\left(a_{0, \infty}-a_{i, j}\right)\right)$ holds for every $j<\omega$, $d^{\star}$ satisfies $\left\{\phi_{i}\left(\operatorname{rv}_{\delta_{n}}\left(x-a_{i, j}\right), b_{i, j}\right)\right\}_{j<\omega}$, which is a contradiction. All rows are inconsistent, which concludes our proof.

With minor modifications, the proof goes through in the case of infinite burden $\lambda$. However, one must be careful regarding the precise statement of this generalisation. Assume we are in mixed characteristic $(0, p)$, and the burden $\lambda$ is of cofinality $\operatorname{cf}(\lambda)=\omega$. Then the very first argument of the proof is no longer true: one cannot necessary assume that there are $\lambda$-many formulas $\tilde{\phi}_{i}\left(x, y_{i}\right)$ in the inp-pattern of the form

$$
\tilde{\phi}_{i}\left(x, c_{i, j}\right)=\phi_{i}\left(\operatorname{rv}_{\delta_{n}}\left(x-a_{i, j ; 1}\right), \ldots, \operatorname{rv}_{\delta_{n}}\left(x-a_{i, j ; k}\right) ; \mathbf{b}_{i, j}\right),
$$

for a certain $n<\omega$. This depends of course of the cofinality of $\lambda$. Nonetheless, this is the only problem. One gets the following statement:

Corollary 3.3. Let $\lambda$ be an infinite cardinal in Card*.

- Let $K$ be a mixed characteristic Henselian valued field. Assume that the union of sorts $\mathrm{RV}_{<\omega}$ with the induced structure is of burden $\lambda$. Then, the field $K$ is of burden $\lambda$ if $\operatorname{cf}(\lambda)>\omega$, and of burden $\lambda$ or $\operatorname{act}(\lambda)$ if $\operatorname{cf}(\lambda)=\omega$.
- Let $K$ be a benign Henselian valued field. Assume that the sort RV with the induced structure is of burden $\lambda$. Then the field $K$ is of burden $\lambda$.

Proof. We treat the case of mixed characteristic Henselian valued field. We prove similarly the case of benign Henselian valued fields. Let $\kappa \geq \lambda$ be the burden of $K$, and let

$$
\left\{\tilde{\phi}_{i}\left(x, y_{i}\right),\left(c_{i, j}\right)_{j<\omega}\right\}_{i<\kappa}
$$

be an inp-pattern of depth $\kappa$. If $\kappa$ is of cofinality $\operatorname{cf}(\kappa)>\omega$, then there are $\kappa$-many formulas $\tilde{\phi}_{i}\left(x, y_{i}\right)$ in the inp-pattern of the form

$$
\tilde{\phi}_{i}\left(x, c_{i, j}\right)=\phi_{i}\left(\operatorname{rv}_{\delta_{n}}\left(x-a_{i, j ; 1}\right), \ldots, \operatorname{rv}_{\delta_{n}}\left(x-a_{i, j ; k}\right) ; \mathbf{b}_{i, j}\right)
$$

for a certain $n<\omega$. We deduce an inp-pattern of depth $\kappa$ in the $\mathrm{RV}_{<\omega}$-sort. Indeed, we follow the exact same proof with few changes in Claim 6:

- The minimum of $\left\{\gamma_{k}\right\}_{k<\lambda}$ may not exist, but one can pick $\gamma$ in an extension of the monster model, realising the cut $\left\{\gamma \in \Gamma \mid \gamma<\gamma_{k}\right.$ for all $\left.k<\lambda\right\} \cup\left\{\gamma \in \Gamma \mid \gamma>\gamma_{k}\right.$ for some $\left.k<\lambda\right\}$. By Claim 5, we have for all $a \in K, \operatorname{val}(a)>\gamma$ implies $\operatorname{val}(a)+\delta_{n}>\gamma_{k}$ for all $k<\lambda$.
- Case A stays the same.
- Case B is slightly different, since an $i$ such that for all $k,\left(a_{k, l}\right)_{l<\omega} \Rightarrow a_{i, \infty}$ or $\left(a_{k,-l}\right)_{l<\omega} \Rightarrow a_{i, \infty}$ does not necessarily exist either. We may distinguish three subcases:

1. there is $i$ such that for $\lambda$-many $k,\left(a_{k, l}\right)_{l<\omega}$ or $\left(a_{k,-l}\right)_{l<\omega}$ pseudo-converge to $a_{i, 0}$. We conclude as in the proof.
2. there is $i<\lambda$ such that $\left(a_{i, j}\right)_{j<\omega}$ pseudo-converges to $a_{k, 0}$ for $\lambda$-many $k$. For such a $k$, we have $\operatorname{val}\left(a_{k, 0}-a_{i, \infty}\right)>\operatorname{val}\left(a_{k, 0}-a_{i, 0}\right) \geq \gamma$, and thus $\operatorname{val}\left(a_{k, 0}-a_{i, \infty}\right)+\delta_{n}>\gamma_{k}$. We may conclude as well.
3. there is $i<\lambda$ such that $\left(a_{i,-j}\right)_{j<\omega}$ pseudo-converges to $a_{k, 0}$ for $\lambda$-many $k$. This is an analogue to Subcase (2) just above, where $a_{i,-\infty}$ is taking the place of $a_{i, \infty}$.

Hence, we get $\lambda=\kappa$.
If $\kappa$ is of cofinality $\omega$, let $\left(\lambda_{k}\right)_{k \in \omega}$ be a sequence of successor cardinals cofinal in $\kappa$. By the previous discussion, we find an inp-pattern in $\mathrm{RV}_{<\omega}$ of depth $\lambda_{k}$ for each $\lambda_{k}$. Hence, $\lambda_{k} \leq \lambda$ and $\kappa=\lambda$ or $\kappa=\operatorname{act}(\lambda)$.

Remark 3.4. - Consider now an enriched Henselian valued field $\mathcal{K}=\left(K, \mathrm{RV}_{<\omega}, \ldots\right)$ of characteristic $(0, p), p \geq 0$ in an $\mathrm{RV}_{<\omega}$-enrichment $\mathrm{L}_{\mathrm{RV}_{<\omega, e}}$ of $\mathrm{L}_{\mathrm{RV}_{<\omega}}$. Then, the above proof still holds. The burden of $K$ is equal (modulo the same subtleties when we consider the burden in Card ${ }^{\star}$ ) to the burden of $\mathrm{RV}_{<\omega} \cup \Sigma_{e}$ with the induced structure, where $\Sigma_{e}$ is the set of new sorts in $\mathrm{L}_{\mathrm{RV}_{<\omega}, e} \backslash \mathrm{~L}_{\mathrm{RV}_{<\omega}}$.

- Similarly, an RV-enriched benign Henselian valued field has the same burden as RV $\cup \Sigma_{e}$ where $\Sigma_{e}$ is the set of new sorts in $\mathrm{L}_{\mathrm{RV}, e} \backslash \mathrm{~L}_{\mathrm{RV}}$.


### 3.1.2 Applications to $p$-adic fields

In this paragraph, $p$ is a prime number. We will deduce from Theorem 3.2, as an application, that any finite extension of $\mathbb{Q}_{p}$ is dp-minimal. This is already known (in fact, all local fields of characteristic 0 are dp-minimal). One can refer to the classification on dp-minimal fields by Will Johnson [22]. The fact that $\mathbb{Q}_{p}$ is dp-minimal is due to Dolich, Goodrick and Lippel [14, Section 6] and Aschenbrenner, Dolich, Haskell, Macpherson and Starchenko in [4, Corollary 7.9.]. In Section 3.3, we will study more generally unramified mixed characteristic Henselian valued fields.

Theorem 3.5. The theory of any finite extension of $\mathbb{Q}_{p}$ in the language of rings is dp-minimal.

A characterisation of dp-minimality is the following: for any mutually indiscernible sequences $\left(a_{i}\right)_{i<\omega}$ and $\left(b_{i}\right)_{i<\omega}$ and any point $c$, one of these two sequences is indiscernible over $c$. As we already mention earlier, a theory is dp-minimal if and only if it is NIP and inp-minimal (see [34, Lemma 1.4] ). Since finite extensions of $\mathbb{Q}_{p}$ are NIP, we have to prove that they are inp-minimal. Recall first that the valuation in a finite extension of $\mathbb{Q}_{p}$ is definable in the language of rings:

$$
\operatorname{val}(x) \geq 0 \Leftrightarrow \exists y 1+\pi x^{q}=y^{q},
$$

where $\pi$ is an element of minimal positive valuation and $q$ is a prime with $q \neq p$. We can safely consider $\mathbb{Q}_{p}$ in the two-sorted language of valued fields $\mathrm{L}=\mathrm{L}_{\mathrm{Mac}} \cup \mathrm{L}_{\text {Pres }} \cup\{$ val $\}$, where $\mathrm{L}_{\mathrm{Mac}}=\mathrm{L}_{\text {Rings }} \cup\left\{P_{n}\right\}_{n \geq 2}$ is the language of Macintyre with a predicate $P_{n}$ for the subgroup of $n$ th-power of $\mathbb{Q}_{p}$ and where LPres is the language of Presburger arithmetic. We have the following well known result, that we already discussed in the example below Proposition 1.10:
Fact 3.6. The theory $\operatorname{Th}\left(\mathbb{Q}_{p}\right)$ eliminates quantifiers. In particular, the value group is a pure sort.
Let $\mathcal{K}=(K, \Gamma)$ be a finite extension of $\mathbb{Q}_{p}$ and let $\pi \in K$ be an element of minimal positive valuation. By interpretability, we obtain:
Remark 3.7. The value group $\Gamma$ is purely stably embedded in $\mathcal{K}$. Since $\Gamma$ is a Z-group (as a finite extension of a $Z$-group), it is in particular inp-minimal.

Fix some $n \in \mathbb{N}$. We have the following exact sequence

$$
1 \longrightarrow \mathcal{O}^{\times} /\left(1+\mathfrak{m}_{\delta_{n}}\right) \longrightarrow \mathrm{RV}_{\delta_{n}}^{\star_{n}} \stackrel{\mathrm{val}_{\mathrm{r}} \delta_{\delta_{n}}}{\longrightarrow} \Gamma \longrightarrow 0,
$$

where $\delta_{n}=\operatorname{val}\left(p^{n}\right)$ and $\mathfrak{m}_{\delta_{n}}=\left\{x \in K \mid \operatorname{val}(x)>\operatorname{val}\left(p^{n}\right)\right\}$. One sees that $\left(\mathcal{O} / \mathfrak{m}_{\delta_{n}}\right)^{\times} \simeq \mathcal{O}^{\times} /\left(1+\mathfrak{m}_{\delta_{n}}\right)$ is finite, or in other words, that the valuation map val $_{\mathrm{rv}_{\delta_{n}}}$ is finite to one. It follows by Lemma 1.26 that $\mathrm{RV}_{\delta_{n}}$ is also inp-minimal. Since this holds for arbitrary $n \in \mathbb{N}$, $R V=\bigcup_{n} R V_{\delta_{n}}$ is inp-minimal. We conclude by using Theorem 3.2.

The next application is a anticipation of the next paragraph. We provide a new proof of the nonuniform definability of an angular component. It can in fact already be deduced from [11]. Recall that an angular component is a group homomorphism ac : $\left(K^{\star}, \cdot\right) \rightarrow\left(k^{\star}, \cdot\right)$ such that $a c_{\mathcal{O}^{\circ} \times}=$ res ${ }_{1 \mathcal{O}} \times$.


Consider any theory $T_{\mathrm{ac}}$ of a valued field endowed with an ac-map, and assume that both the value group $\Gamma$ and residue field $k$ are infinite. Then by Fact 1.25 and bi-interpretability on unary sets, on sees that the RV-sort is of burden at least 2. The set $\mathrm{RV}^{\star}$ is indeed in definable bijection with the direct product $\Gamma \times k^{\star}$.

In the field of $p$-adics $\mathbb{Q}_{p}$, an angular component ac is definable in the language of rings. We can show now easily that this definition cannot be uniform:
Corollary 3.8. There is no formula which gives a uniform definition of an ac-map in $\mathbb{Q}_{p}$ for every prime $p$.

Notice that this has already been observed by Pas in [28].
Proof. By Chernikov-Simon [11], we know that the ultraproduct of $p$-adic $\mathcal{F}=\prod_{\mathcal{U}} \mathbb{Q}_{p}$, where $\mathcal{U} \subset \mathcal{P}$ is an ultrafilter on the set of primes, is inp-minimal in the language of rings (recall that the p-adic valuation is uniformly definable in $\mathrm{L}_{\text {Rings }}$ ). The residue field and the value group are infinite since they are respectively a pseudo-finite field and a $\mathbb{Z}$-group. By the above discussion, the ac-map cannot be defined in the language of rings, as it would contradict inp-minimality.

### 3.2 Benign Henselian Valued Fields

Let $\mathcal{K}=(K, \Gamma, k)$ be a saturated enough benign Henselian valued field. We will compute the burden of $\mathrm{RV}:=K^{\star} / 1+\mathfrak{m}$ in terms of the burden of $k$ and $\Gamma$. As the RV-sort is stably embedded, we will consider it as a structure on its own. By Fact 1.47, the induced structure is given by:

$$
\left\{\mathrm{RV},(k, \cdot,+, 0,1),(\Gamma,+, 0,<), \mathrm{val}_{\mathrm{rv}}: \mathrm{RV} \rightarrow \Gamma, k^{\star} \rightarrow \mathrm{RV}\right\} .
$$

Notice in particular that there is no need of the symbol $\oplus$ as we consider the sort $k$ and $\Gamma$ instead. The language is denoted by L. In other words the sort RV is no more than an enriched exact sequence of abelian groups:

$$
1 \rightarrow k^{\star} \rightarrow \mathrm{RV}^{\star} \xrightarrow{\text { val }_{\mathrm{vv}}} \Gamma \rightarrow 0,
$$

where $k=k^{\star} \cup\{0\}$ is endowed with its field structure and $\Gamma$ is endowed with its ordered abelian group structure. As $\Gamma$ is torsion free, $k^{\star}$ is a pure subgroup of $\mathrm{RV}^{\star}$. The idea to consider RV as an enrichment of abelian groups is already present in [11] and has been developed in [3].

### 3.2.1 Reduction of burden to $\Gamma$ and $k$

Let us recall a result of Chernikov and Simon:
Theorem 3.9 ([11, Theorem 1.4]). Assume $\mathcal{K}$ is a Henselian valued field of equicharacteristic 0. Assume the residue field $k$ satisfies

$$
\begin{equation*}
k^{\star} /\left(k^{\star}\right)^{p} \text { is finite for every prime } p . \tag{k}
\end{equation*}
$$

Then $\mathcal{K}$ is inp-minimal if and only if RV with the induced structure is inp-minimal if and only if $k$ and $\Gamma$ are both inp-minimal.

It will now be easy to extend this theorem. We have already seen the reduction to the RV-sort for any benign Henselian valued field, without the assumption $\left(H_{k}\right)$. For the reduction to $\Gamma$ and $k$, one can give first an easy bound, also independent of the assumption $\left(H_{k}\right)$. Indeed, recall that in an $\aleph_{1}$-saturated model, any pure exact sequence of abelian groups splits (Fact 1.73). In particular, there exists a section $\mathrm{ac}_{\mathrm{rv}}: \mathrm{RV}^{\star} \rightarrow k^{\star}$ of the valuation val ${ }_{\mathrm{rv}}$ or equivalently, there exists an angular component ac : $K^{\star} \rightarrow k^{\star}$ (as we already discussed in Paragraph 3.1.2).


Recall that $\mathrm{L}_{\mathrm{ac}}$ is the language L extended by a unary function $\mathrm{ac}_{\mathrm{rv}}: k \rightarrow \mathrm{RV}$. A direct translation of Fact 2.1 gives:

Fact 3.10 (Trivial bound). We have $\operatorname{bdn}_{\mathrm{L}}(\Gamma)=\operatorname{bdn}_{\mathrm{Lac}_{\mathrm{ac}}}(\Gamma)$ and $\operatorname{bdn}_{\mathrm{L}}(k)=\operatorname{bdn}_{\mathrm{Lac}_{\mathrm{a}}}(k)$ as well as the following:

$$
\operatorname{bdn}_{\mathrm{L}}(\mathrm{RV}) \leq \operatorname{bdn}_{\mathrm{L}_{\mathrm{ac}}}(\mathrm{RV})=\operatorname{bdn}_{\mathrm{L}}(\Gamma)+\operatorname{bdn}_{\mathrm{L}}(k) .
$$

A valued field $\mathcal{K}_{\text {ac }}$ together with an angular component ac can be considered as an RV-enrichment of $\mathcal{K}$. Using the enriched version of Theorem 3.2 (see Remark 3.4), we get:

Theorem 3.11. Let $\mathcal{K}_{\mathrm{ac}}=(\mathcal{K}, \Gamma, k, \mathrm{val}, \mathrm{ac})$ be a benign Henselian valued field endowed with an ac-map. Then:

$$
\operatorname{bdn}\left(\mathcal{K}_{\mathrm{ac}}\right)=\operatorname{bdn}(k)+\operatorname{bdn}(\Gamma) .
$$

If we do not want to consider an ac-map, we can also compute the burden using the torsion-free case of Theorem 2.2 together with Theorem 3.2. We get:

Theorem 3.12. Let $\mathcal{K}$ be a benign Henselian valued field. Then:

$$
\operatorname{bdn}(\mathcal{K})=\max _{n \geq 0}\left(\operatorname{bdn}\left(k^{\star} / k^{\star n}\right)+\operatorname{bdn}(n \Gamma)\right) .
$$

Using Adler's convention, we have the same formula:

$$
\operatorname{bdn}^{\star}(\mathcal{K})=\max _{n \geq 0}\left(\operatorname{bdn}^{\star}\left(k^{\star} / k^{\star n}\right)+\operatorname{bdn}^{\star}(n \Gamma)\right) .
$$

This gives a full answer to [11, Problem 4.3] and [11, Problem 4.4]:
Corollary 3.13. Let $\mathcal{K}=(K, \mathrm{RV}, k, \Gamma)$ be a benign Henselian valued field. Assume that:

$$
\begin{equation*}
k^{\star} /\left(k^{\star}\right)^{p} \text { is finite for every prime } \mathrm{p} . \tag{k}
\end{equation*}
$$

Then we have the equalities

$$
\operatorname{bdn}(\mathcal{K})=\operatorname{bdn}(\operatorname{RV})=\max (\operatorname{bdn}(k), \operatorname{bdn}(\Gamma)) .
$$

Also, in the case that $\mathcal{K}$ is not trivially valued, the value group $\Gamma$ is necessary of burden $\operatorname{bdn}(\Gamma)>0$. It follows that a non-trivially valued benign Henselian field $\mathcal{K}$ is inp-minimal if and only if $\Gamma, k$ are inp-minimal and $k$ satisfies $\left(H_{k}\right)$.

Similarly to the proof of non-existence of a uniform definition of the angular component of $\mathbb{Q}_{p}$, we can notice the following:

Remark 3.14. Let $\mathcal{K}$ be a benign Henselian valued field of finite burden. Assume that the residue field is infinite and satisfies $\left(H_{k}\right)$. Then, no angular component is definable in the language of valued fields $\mathrm{L}_{d i v}$.

The reason is of course that in such a case, the two terms $\max (\operatorname{bdn}(k), \operatorname{bdn}(\Gamma))$ and $\operatorname{bdn}(k)+\operatorname{bdn}(\Gamma)$ are distinct.

All these results hold resplendently. In fact by definition, a benign Henselian valued field can have an enriched value group and residue field. Let us clarify by stating the previous theorem in an enriched language:

Remark 3.15. If $\mathcal{K}=(K, \operatorname{RV}, k, \Gamma, \ldots)$ is a $\{\Gamma\}-\{k\}$-enriched benign Henselian valued field in a $\{\Gamma\}$ - $\{k\}$-enrichment $\mathrm{L}_{\Gamma, k, e}$ of $\mathrm{L}_{\Gamma, k}$, then

$$
\operatorname{bdn}(\mathcal{K})=\operatorname{bdn}\left(\operatorname{RV} \cup \Sigma_{e}\right)=\max _{n}\left(\operatorname{bdn}\left(k^{\star} / k^{\star n}\right)+\operatorname{bdn}(n \Gamma), \operatorname{bdn}\left(\Sigma_{e}\right)\right),
$$

where $\Sigma_{e}$ is the set of new sorts in $\mathrm{L}_{\Gamma, k, e} \backslash \mathrm{~L}_{\Gamma, k}$.
To conclude this short subsection, let us discuss more on the hypothesis $\left(H_{k}\right)$ and bounded fields.

### 3.2.2 Bounded fields and applications

A bounded field is a field with finitely many extensions of degree $n$ for every integer $n$. The absolute Galois group is called small if it contains finitely many open subgroups of index $n$. These two conditions are equivalent for perfect fields: a perfect field is bounded if and only if its absolute Galois group is small. Such a field $K$ satisfies in particular the following:

$$
\begin{equation*}
K^{\star} /\left(K^{\star}\right)^{p} \text { is finite for every prime } \mathrm{p}, \tag{H}
\end{equation*}
$$

(see for example [16, Proposition 2.3]), and it's clear that $(H)$ implies $\left(H_{k}\right)$. It also implies:

$$
\Gamma / p \Gamma \text { is finite for every prime } \mathrm{p},
$$

The condition $H_{k}$ might be restrictive but it allows various burdens for the residue field. However, the condition $H_{\Gamma}$ implies inp-minimality for the value group. Indeed, an abelian group $\Gamma$ satisfying $\left(H_{\Gamma}\right)$ is called non-singular. In the pure structure of ordered abelian groups, non-singular ordered abelian groups are exactly the dp-minimal ones (see [21, Theorem 5.1]). We have the following examples :
Examples. - The Hahn field $\mathbb{F}_{p}^{a l g}((\mathbb{Z}[1 / p]))$ is algebraically maximal Kaplansky Henselian. By Jahnke, Simon and Walsberg, the value group $\mathbb{Z}[1 / p]$ is inp-minimal as it satisfies $\left(H_{\Gamma}\right)$. The residue field $\mathbb{F}_{p}^{\text {alg }}$ satisfies $\left(H_{k}\right)$ and is inp-minimal. By Theorem 3.12, this Hahn field is inpminimal.

- In general, a bounded benign Henselian valued field $\mathcal{K}$ with residue field $k$ has burden $\max (\operatorname{bdn}(k), 1)$.
Montenegro has computed the burden of some theories of bounded fields, namely bounded pseudo real closed fields (PRC fields) and pseudo $p$-adicaly closed fields (PpC fields). We recall here these theorems (see [27, Theorems $4.22 \& 4.23]$ ):
Theorem 3.16. Let $k$ be a bounded PRC field. Then $\operatorname{Th}(k)$ is $\mathrm{NTP}_{2}$, strong and of burden the (finite) number of orders in $k$.
Theorem 3.17. Let $k$ be an PpC field. Then $\operatorname{Th}(k)$ is $\mathrm{NTP}_{2}$ if and only if $\operatorname{Th}(k)$ is strong if and only if $k$ is bounded. In this case, the burden of $\operatorname{Th}(k)$ is the (finite) number of $p$-adic valuations in $k$.


### 3.3 Unramified mixed characteristic Henselian Valued Fields

Let $\mathcal{K}=\left(K, \mathrm{RV}_{<\omega}, \Gamma, k\right)$ be an unramified Henselian valued field of characteristic ( $0, p$ ), $p \geq 2$ with perfect residue field $k$. We denote by 1 the valuation of $p$. The value group $\Gamma$ contains $\mathbb{Z} \cdot 1$ as a convex subgroup. Recall that in this context, it is more convenient to denote the $n^{\text {th }} \mathrm{RV}$-sort by $\mathrm{RV}_{n}^{\star}:=K^{\star} /\left(1+\mathfrak{m}^{n}\right)$ where $\mathfrak{m}=\{x \in K \mid \operatorname{val}(x)>0\}$ is the maximal ideal of the valuation ring $\mathcal{O}$. Notice that $\mathfrak{m}^{n}=p^{n} \mathcal{O}$ for every integer $n$. Similarly to the previous section, we will compute the burden of $\mathrm{RV}_{<\omega}=\cup_{n \in \mathbb{N}} \mathrm{RV}_{n}$ in terms of the burden of $k$ and $\Gamma$.

### 3.3.1 Reduction from $\mathrm{RV}_{<\omega}$ to $\Gamma$ and $k$

Now we can look for the burden of $\mathrm{RV}_{n}$. We start with a harmless observation:
Observation 3.18. Let $m<n$ be integers. The element $p^{m}$ is of valuation $m$. By [17, Proposition 2.8], $\mathrm{RV}_{m}$ is $\emptyset$-interpretable in $\mathrm{RV}_{n}$, with base set $\mathrm{RV}_{n}$ quotiented by an equivalence relation. Hence the burden of $\mathrm{RV}_{n}$ can only grow with $n$ : for $m<n, \operatorname{bdn}\left(\mathrm{RV}_{m}\right) \leq \operatorname{bdn}\left(\mathrm{RV}_{n}\right)$.

Recall that in this context of unramified mixed characteristic Henselian valued fields with perfect residue field, the $n^{\text {th }}$ residue ring $\mathcal{O}_{n}:=\mathcal{O} / p^{n} \mathcal{O}$ is isomorphic to the $n$-truncated ring of Witt vectors (see Proposition 1.67). We work now in the following languages:

$$
\begin{aligned}
\mathrm{L}=\{ & K, \Gamma,\left(\mathrm{RV}_{n}\right)_{n<\omega},\left(W_{n}(k)\right)_{n<\omega}, \mathrm{val}: K^{\star} \rightarrow \Gamma \\
& \left.\left(\operatorname{res}_{n}: \mathcal{O} \rightarrow W_{n}(k)\right)_{n<\omega},\left(\operatorname{rv}_{n}: K^{\star} \rightarrow \operatorname{RV}_{n}\right)_{n<\omega}\right\},
\end{aligned}
$$

which is a little variation of (and bi-interpretable with) the language $\mathrm{L}_{\mathrm{RV}}{ }_{<\omega}$, where the structure of the $\mathrm{RV}_{n}$ 's is described with exact sequences. We can also add the ac-maps to this language:

$$
\mathrm{L}_{a c_{<\omega}}=\mathrm{L} \cup\left\{\left(\mathrm{ac}_{n}: K^{\star} \rightarrow W_{n}(k)\right)_{n<\omega}\right\}
$$

Here is a consequence of Corollary 1.62, Remark 1.16 and Fact 1.18:

Corollary 3.19. We have:

- $\operatorname{bdn}\left(W_{n}(k)\right)=\kappa_{\text {inp }}^{1}\left(W_{n}(k)\right)=\kappa_{\text {inp }}^{n}(k)$.
- $\operatorname{bdn}((W(k),+, \cdot, \pi: W(k) \rightarrow k))=\kappa_{\text {inp }}^{\aleph_{0}}(k)$.

Recall that we have the following inequalities (see Paragraph 1.1.2):

$$
n \cdot \kappa_{\text {inp }}^{1}(k) \leq \kappa_{\text {inp }}^{n}(k) \quad \kappa_{\text {inp }}^{n}(k)+1 \leq\left(\kappa_{\text {inp }}^{1}(k)+1\right)^{n}
$$

In particular, if $k$ is infinite then the burden of $\left(W_{n}(k),+, \cdot, \pi\right)$ is at least $n$.
In the language $\mathrm{L}_{\mathrm{ac}<\omega}$, a consequence of Proposition 1.67 is that, for every $n<\omega$ the sort $W_{n}(k)$ is pure (in particular stably embedded) and orthogonal to $\Gamma$, as it is $\emptyset$-bi-interpretable with $\left(k^{n},+, \cdot, p_{i}, i<\right.$ $n$ ), which is a pure sort orthogonal to $\Gamma$. It follows that $W_{n}(k)$ doesn't have more structure in $\mathrm{L}_{\mathrm{ac}<\omega}$ than in L. Similarly, the burden of $\Gamma$ is the same in any of the above languages. Hence, we actually have the following equalities:

$$
\begin{align*}
\operatorname{bdn}_{\mathrm{L}}\left(W_{n}(k)\right) & =\operatorname{bdn}_{\mathrm{L}_{\mathrm{ac}<\omega}}\left(W_{n}(k)\right)  \tag{7}\\
\operatorname{bdn}_{\mathrm{L}}(\Gamma) & =\operatorname{bdn}_{\mathrm{Lac}_{\mathrm{a}}<\omega}(\Gamma) \tag{8}
\end{align*}
$$

We are now able to give a relationship between $\operatorname{bdn}\left(\operatorname{RV}_{n}\right)$ and $\operatorname{bdn}\left(W_{n}(k)\right)$.
Proposition 3.20. [Trivial bound] We have

$$
\begin{aligned}
\max \left(\operatorname{bdn}_{\mathrm{L}}\left(W_{n}(k)\right), \operatorname{bdn}_{\mathrm{L}}(\Gamma)\right) \leq \operatorname{bdn}_{\mathrm{L}}\left(\operatorname{RV}_{n}\right) & \leq \operatorname{bdn}_{\mathrm{Lac}_{<\omega}}\left(\operatorname{RV}_{n}\right) \\
& =\operatorname{bdn}_{\mathrm{L}}\left(W_{n}(k)\right)+\operatorname{bdn}_{\mathrm{L}}(\Gamma)
\end{aligned}
$$

Proof. By Proposition 1.67, we have the exact sequence of abelian groups:

$$
1 \rightarrow W_{n}(k)^{\times} \rightarrow \mathrm{RV}_{n}^{\star} \rightarrow \Gamma \rightarrow 0
$$

The first inequality is clear if one shows that $\operatorname{bdn}_{\mathrm{L}}\left(W_{n}(k)\right)=\operatorname{bdn}_{\mathrm{L}}\left(W_{n}(k)^{\times}\right)$where $W_{n}(k)^{\times}$is endowed with the induced structure. The second inequality is also clear, as adding structure can only make the burden grow. Let $\left\{\phi_{i}\left(x, y_{i}\right),\left(a_{i, j}\right)_{j<\omega}\right\}_{i \in \lambda}$ be an inp-pattern in $W_{n}(k)$, with $\left(a_{i, j}\right)_{i<\lambda, j<\omega}$ mutually indiscernible. Let $d \models\left\{\phi\left(x, a_{i, 0}\right)\right\}_{i \in \lambda}$ be a realisation of the first column. In the case where $d \in W_{n}(k)^{\times}$, there is nothing to do. Otherwise, $1+d \in W_{n}(k)^{\times}$and $\left\{\phi_{i}\left(x-1, y_{i}\right),\left(a_{i, j}\right)_{j<\omega}\right\}_{i \in \lambda}$ is an inp-pattern in $W_{n}(k)^{\times}$of depth $\lambda$. This concludes the proof of the first inequality.

We work now in $\mathrm{L}_{\mathrm{ac}<\omega}$, where we interpret $\left(a c_{n}\right)_{n}$ as a compatible sequence of angular components (it exists by $\aleph_{1}$-saturation). Recall that the burden may only increase. Then, the above exact sequences (definably) split in $\mathrm{L}_{\mathrm{ac}<\omega}$, as we add a section. By the previous discussion, $W_{n}(k)^{\times}$and $\Gamma$ are orthogonal and stably embedded. We apply now Fact 1.25 : the burden $\operatorname{bdn}_{\mathrm{Lac}^{2}<\omega}\left(\mathrm{RV}_{n}^{\star}\right)$ is equal to $\operatorname{bdn}_{\mathrm{Lac}_{<\omega}}\left(W_{n}(k)^{\times}\right)+\operatorname{bdn}_{\mathrm{L}_{\mathrm{ac}<\omega}}(\Gamma)=\operatorname{bdn}_{\mathrm{L}}\left(W_{n}(k)\right)+\operatorname{bdn}_{\mathrm{L}}(\Gamma)$.

Combining Corollary 3.3, Corollary 3.19 and Proposition 3.20 , one gets:
Theorem 3.21. Let $\mathcal{K}=(K, k, \Gamma)$ be an unramified mixed characteristic Henselian valued field. We denote by $\mathcal{K}_{\mathrm{ac}<\omega}=\left(K, k, \Gamma, \mathrm{ac}_{n}, n<\omega\right)$ the structure $\mathcal{K}$ endowed with compatible ac-maps. Assume the residue field $k$ is perfect. One has

$$
\operatorname{bdn}(\mathcal{K})=\operatorname{bdn}\left(\mathcal{K}_{\mathrm{ac}<\omega}\right)=\max \left(\aleph_{0} \cdot \operatorname{bdn}(k), \operatorname{bdn}(\Gamma)\right)
$$

And its enriched version:

Remark 3.22. Let $\mathrm{L}_{e}$ be a $\{\Gamma\}$ - $\{k\}$-enrichment of L . Let $\mathcal{K}=(K, k, \Gamma, \ldots)$ be an enriched unramified mixed characteristic Henselian valued field in the language $\mathrm{L}_{e}$. Assume the residue field $k$ is perfect. We denote by $\mathcal{K}_{\mathrm{ac}_{<\omega}}=\left(K, k, \Gamma, \mathrm{ac}_{n}, n<\omega, \ldots\right)$ the structure $\mathcal{K}$ endowed with compatible ac-maps. One has

$$
\operatorname{bdn}(\mathcal{K})=\operatorname{bdn}\left(\mathcal{K}_{\mathrm{ac}<\omega}\right)=\max \left(\aleph_{0} \cdot \operatorname{bdn}(k), \operatorname{bdn}(\Gamma), \operatorname{bdn}\left(\Sigma_{e}\right)\right),
$$

where $\Sigma_{e}$ is the set of new sorts in $\mathrm{L}_{e} \backslash \mathrm{~L}$.
This is a simple calculation, unless we want to consider burden in Card ${ }^{\star}$.
Remark 3.23. Let $\mathcal{K}$ and $\mathcal{K}_{\mathrm{ac}<\omega}$ as above. We consider the second definition of burden (Definition 1.29). We have

$$
\operatorname{bdn}^{\star}(\mathcal{K})=\operatorname{bdn}^{\star}\left(\mathcal{K}_{\mathrm{ac}<\omega}\right)=\max \left(\aleph_{0} \cdot{ }^{\star} \operatorname{bdn}^{\star}(k), \operatorname{bdn}^{\star}(\Gamma)\right) .
$$

It follows that an unramified mixed characteristic valued field with an infinite perfect residue field is never strong.

This is what we will prove now.
Proof. We use the same notation as before in this section. Unfortunately, due to the ambiguity in Corollary 3.3 concerning $\operatorname{bdn}^{\star}(\mathcal{K}) \in\left\{\operatorname{bdn}^{\star}\left(\operatorname{RV}_{<\omega}\right)\right.$, act $\left.\left(\operatorname{bdn}^{\star}\left(\operatorname{RV}_{<\omega}\right)\right)\right\}$ in the case that $\operatorname{cf}\left(\operatorname{bdn}\left(\operatorname{RV}_{<\omega}\right)\right)=\omega$, we have to go back to the proof of Theorem 3.2.

We first show that $\operatorname{bdn}^{\star}(\mathcal{K})$ is at least $\aleph_{0} \bullet^{\star} \operatorname{bdn}^{\star}(k)$. Recall that $W_{n}(k) \simeq \mathcal{O}_{n}:=\mathcal{O} / \mathfrak{m}^{n}$ is interpretable (with one-dimensional base set $\mathcal{O} \subset K$ ), and so is the projective system $\left\{W_{n}(k), \pi_{n, m}\right.$ : $\left.W_{n}(k) \rightarrow W_{m}(k), n>m\right\}$ and the projection maps $\chi_{n, n}: W_{n}(k) \rightarrow k, x=\left(x_{1}, \ldots, x_{n}\right) \mapsto x_{n}$. If $\operatorname{cf}(\operatorname{bdn}(k))>\aleph_{0}$, there is nothing to do as $\operatorname{bdn}^{\star}(k)=\aleph_{0}{ }^{\star} \operatorname{bdn}^{\star}(k)$. Assume $\operatorname{cf}(\operatorname{bdn}(k)) \leq \aleph_{0}$. We write $\operatorname{bdn}^{\star}(k)=\sup _{n<\omega} \lambda_{n}$ with $\lambda_{n} \in$ Card. Let $P_{n}\left(x_{k}\right)$ be an inp-pattern with $x_{k} \in k,\left|x_{k}\right|=1$, of depth $\lambda_{n}$ for every $n \in \omega$. Then, the pattern $P(x)=\cup_{n \in \omega} P_{n}\left(\chi_{n, n}\left(\pi_{n}(x)\right)\right)$ is an inp-pattern in $K$ of depth $\aleph_{0}{ }^{\star} \operatorname{bdn}^{\star}(k)$. One gets:

$$
\operatorname{bdn}^{\star}(\mathcal{K}) \geq \max \left(\aleph_{0} \cdot{ }^{\star} \operatorname{bdn}^{\star}(k), \operatorname{bdn}^{\star}(\Gamma)\right) .
$$

We now prove that $\max \left(\aleph_{0} .^{\star} \operatorname{bdn}^{\star}(k), \operatorname{bdn}^{\star}(\Gamma)\right)$ is an upper bound for $\operatorname{bdn}^{\star}\left(\mathcal{K}_{\mathrm{ac}<\omega}\right)$.
Case 1: $\aleph_{0} .^{\star} \operatorname{bdn}^{\star}(k) \geq \operatorname{bdn}^{\star}(\Gamma)$.
Subcase 1.A: $\operatorname{cf}(\operatorname{bdn}(k))>\aleph_{0}$. By Corollary 3.3, $\operatorname{bdn}^{\star}\left(\mathcal{K}_{\mathrm{ac}<\omega}\right)=\operatorname{bdn}^{\star}\left(\mathrm{RV}_{<\omega}\right)=$ $\sup _{n}\left(\kappa_{\text {inp }}^{n}(k), \operatorname{bdn}^{\star}(\Gamma)\right)=\operatorname{bdn}^{\star}(k)=\aleph_{0} \cdot \operatorname{bdn}^{\star}(k)$. We used the submultiplicativity of the burden, which gives here $\kappa_{\text {inp }}^{n \star}(k)=\kappa_{\text {inp }}^{1 \star}(k)=\operatorname{bdn}^{\star}(k)$ for all $n \in \mathbb{N}$.
Subcase 1.B: $\operatorname{cf}(\operatorname{bdn}(k)) \leq \aleph_{0}$. Then $\operatorname{act}\left(\operatorname{bdn}^{\star}\left(\operatorname{RV}_{<\omega}\right)\right)=\aleph_{0} \cdot \operatorname{bdn}(k)$. By Corollary 3.3, we have $\operatorname{bdn}^{\star}\left(\mathcal{K}_{\mathrm{ac}<\omega}\right) \leq \aleph_{0} \cdot \operatorname{bdn}(k)$.
Case 2: $\operatorname{bdn}^{\star}(\Gamma)>\aleph_{0} \cdot^{\star} \operatorname{bdn}^{\star}(k)$. If $\operatorname{bdn}^{\star}(\Gamma)$ is in Card, this is clear by Corollary 3.3. Assume $\operatorname{bdn}^{\star}(\Gamma)$ is of the form $\lambda_{-}$for a limit cardinal $\lambda \in$ Card. Notice that this case occurs only if the sort $\Gamma$ is enriched. We work in the corresponding enrichment of language $\mathrm{L}_{\mathrm{RV}}<\omega$ together with $\mathrm{ac}_{n}$-maps. We have to show that $\lambda_{-}$is an upper bound for $\operatorname{bdn}^{\star}\left(\mathcal{K}_{\text {ac< }}\right)$. Let $P(x)=\left\{\theta_{i}\left(x, y_{i, j}\right),\left(c_{i, j}\right)_{j \in \overline{\mathbb{Z}}}\right\}_{i \in \lambda}$ be an inp-pattern in $K$ of depth $\lambda$ with $|x|=1$ and $\left(c_{i, j}\right)_{i<\lambda, j \in \overline{\mathbb{Z}}}$ be a mutually indiscernible array. Then, by Fact 1.64, one can assume that each formula $\theta_{i}\left(x, c_{i, j}\right)$ in $P(x)(i \leq \lambda, j \in \overline{\mathbb{Z}})$ is of the form

$$
\tilde{\theta}_{i}\left(\operatorname{rv}_{n_{i}}\left(x-\alpha_{i, j}^{1}\right), \ldots, \operatorname{rv}_{n_{i}}\left(x-\alpha_{i, j}^{m}\right), \beta_{i, j}\right),
$$

for some integers $n_{i}$ and $m$, and where $\alpha_{i, j}^{1}, \ldots, \alpha_{i, j}^{m} \in K, \beta_{i, j} \in \mathrm{RV}_{n_{i}}$ and $\tilde{\theta}_{i}$ is an $\mathrm{RV}_{n_{i}}$-formula. As in the proof of Theorem 3.2, we may assume with no restriction that $m=1$. As $\mathrm{RV}_{n_{i}}=W_{n_{i}}(k)^{\times} \times \Gamma$ is the direct product of the orthogonal and stably embedded sorts $W_{n_{i}}(k)^{\times}$and $\Gamma$, we may assume $\theta_{i}\left(x, c_{i, j}\right)$ is equivalent to a formula of the form

$$
\phi_{i}\left(\mathrm{ac}_{n_{i}}\left(x-\alpha_{i, j}\right), a_{i, j}\right) \wedge \psi_{i}\left(\operatorname{val}\left(x-\alpha_{i, j}\right), b_{i, j}\right)
$$

where $\phi_{i}\left(x_{W_{n_{i}}}, a_{i, j}\right)$ is a $W_{n_{i}}$-formula and $\psi_{i}\left(x_{\Gamma}, b_{i, j}\right)$ is a $\Gamma$-formula. By Claim 6 in Theorem 3.2 (or more precisely, by a generalisation of Claim 6 to infinite depth $M=\lambda$ ), one may assume that there is $k<\lambda$ such that for all $i<\lambda$,

$$
\operatorname{val}\left(d-\alpha_{i, 0}\right) \leq \min \left\{\operatorname{val}\left(d-\alpha_{k, \infty}\right), \operatorname{val}\left(\alpha_{k, \infty}-\alpha_{i, 0}\right)\right\}+\max \left(n_{i}, n_{k}\right)
$$

It follows that, if $\operatorname{val}\left(d-\alpha_{k, \infty}\right)=\operatorname{val}\left(\alpha_{k, \infty}-\alpha_{i, 0}\right), \operatorname{val}\left(d-\alpha_{i, 0}\right)$ is equal to $\operatorname{val}\left(d-\alpha_{k, \infty}\right)+n_{i}^{\prime}$ for some $0 \leq n_{i}^{\prime} \leq \max \left(n_{i}, n_{k}\right)$. Otherwise, one has $\operatorname{val}\left(d-\alpha_{i, 0}\right)=\min \left\{\operatorname{val}\left(d-\alpha_{k, \infty}\right), \operatorname{val}\left(\alpha_{k, \infty}-\alpha_{i, 0}\right)\right\}$. We can centralise $P(x)$ in $\alpha_{k, \infty}$, i.e. we can assume that each formula in $P(x)$ is of the form

$$
\phi_{i}\left(\mathrm{ac}_{2 n_{i}}\left(x-\alpha_{k, \infty}\right), a_{i, j}\right) \wedge \psi_{i}\left(\operatorname{val}\left(x-\alpha_{k, \infty}\right), b_{i, j}\right)
$$

(we add new parameters $\operatorname{val}\left(\alpha_{k, \infty}-\alpha_{i, j}\right)$ and $\operatorname{ac}_{2 n_{i}}\left(\alpha_{k, \infty}-\alpha_{i, j}\right)$. Notice that once the difference of the valuation is known, $\operatorname{ac}_{n_{i}}\left(d-\alpha_{i, j}\right)$ can be computed in terms of $\operatorname{ac}_{2 n_{i}}\left(d-\alpha_{k, \infty}\right)$ and $\left.\operatorname{ac}_{2 n_{i}}\left(\alpha_{i, j}-\alpha_{k, \infty}\right)\right)$. By indiscernibility, at least one of the following sets

$$
\left\{\phi_{i}\left(x_{W_{2 n_{i}}}, a_{i, j}\right)\right\}_{j<\omega}
$$

and

$$
\left\{\psi_{i}\left(x_{\Gamma}, b_{i, j}\right)\right\}_{j<\omega}
$$

is inconsistent. Since $\lambda>\sup _{n} \operatorname{bdn}^{\star}\left(W_{n}(k)\right)$, we may assume that

$$
\left.\left\{\psi_{i}\left(x_{\Gamma}, y_{i}\right),\left(b_{i, j}\right)_{j \in \overline{\mathbb{Z}}}\right)\right\}_{i<\lambda}
$$

is an inp-pattern in $\Gamma$. This is a contradiction. Hence, we have $\operatorname{bdn}^{\star}(\mathcal{K})=\lambda_{-}$.
We end now with examples:
Examples. 1. Assume that $k$ is an algebraically closed field of characteristic $p$, and $\Gamma$ is a $\mathbb{Z}$-group. Then $\Gamma$ is inp-minimal, i.e. of burden one (as it is quasi-o-minimal), and one has $\kappa_{i n p}^{n}(k)=n$. By Theorem 3.21, any Henselian mixed characteristic valued field of value group $\Gamma$ and residue field $k$ has burden $\aleph_{0}$. In particular, the quotient field $Q(W(k))$ of the Witt vectors $W(k)$ over $k$ is not strong.
2. Consider once again the field of $p$-adics $\mathbb{Q}_{p}$. We have $\kappa_{\text {inp }}^{n}\left(\mathbb{F}_{p}\right)=0$ for all $n$, and $\operatorname{bdn}(\mathbb{Z})=1$. Then Theorem 3.21 gives $\operatorname{bdn}\left(\mathbb{Q}_{p}\right)=1$.

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[^1]:    ${ }^{1}$ It was later show by Gurevich and Schmitt that any pure ordered abelian group is NIP.

[^2]:    ${ }^{2}$ This happens for instance with the residue field $k$ of an equicharacteristic 0 Henselian valued field: it is closed in the traditional 3-sorted language of valued fields but it is also natural to interpret in addition a short exact sequences of abelian groups $1 \rightarrow k^{\times} \rightarrow \mathrm{RV}^{\star} \rightarrow \Gamma \rightarrow 0$. See Paragraph 1.2.1

[^3]:    ${ }^{3}$ In the literature, the notion of bi-interpretability usually requires furthermore that the composition of the interpretation $f \circ g: \mathcal{M} \rightarrow \mathcal{M}$ (resp. $g \circ f: \mathcal{N} \rightarrow \mathcal{N}$ ) is definable in $\mathcal{M}$ (resp. in $\mathcal{N}$ ). This is however not relevant for the computation of the burden, and we omit it in our definition.

[^4]:    ${ }^{4}$ The Ax-Kochen-Ershov property and relative quantifier elimination for Henselian unramified mixed characteristic valued fields (with possibly imperfect residue field) has been proved in [2].

[^5]:    ${ }^{5}$ The proof shows that one does not need the function $\iota: A \rightarrow B$ in order to describe definable sets in $A$. In a certain sense, $<A>$ is a 'closure' of $A$, as it describes the induced structure on $A$, with no resort to any symbol from $\left.\mathrm{L}_{q} \backslash \mathrm{~L}_{q}\right|_{<A>}$.

