# Multialgebraic structures and applications in abstract theories of quadratic forms and graded rings 

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## Resumo

ROBERTO, K. M. A. Estruturas multi-algébricas e aplicações em teorias abstratas de formas quadráticas e anéis graduados. 2023. 214 f . Tese (Doutorado) - Instituto de Matemática e Estatística, Universidade de São Paulo, São Paulo, 2023.

O objetivo deste trabalho é iniciar a investigação de possíveis relações matemáticas que configurem um "novo quadro adjunto" entre grupos especiais, anéis graduados, 2-grupos e 2 -grupos profinitos. Nós focamos na primeira parte deste programa, i.e. nas relações entre anéis graduados e grupos especiais. Em nossas investigações, a teoria dos multi anéis/hipercorpos desempenhou um papel central, e obtivemos um resultado interessante: uma ampla extensão (para todos os grupos/hipercorpos especiais) da validade do Arason-Pfister Hauptsatz ([7]) - uma resposta positiva para uma questão formulada por J. Milnor no clássico artigo de 1970 ([52]) - e aplicamos este resultado na obtenção de propriedades associadas aos anéis graduados provenientes de hipercorpos especiais ([30], [18]).

Palavras-chave: grupo especial, anéis graduados, 2-grupos profinitos, multi anéis, hipercorpos, formas quadráticas, Arason-Pfister hauptsatz.

## Abstract

ROBERTO, K. M. A. Multialgebraic structures and applications in abstract theories of quadratic forms and graded rings. 2023. 214 f . Tese (Doutorado) - Instituto de Matemática e Estatística, Universidade de São Paulo, São Paulo, 2023.

The aim of this work is to initiate the investigation of the precise mathematical relationships so that possible configuring a "new adjoint" situation between special groups, graded rings, 2 -groups and profinite 2-groups. We focused in the first part of this program, i,e, in the relations between graded rings and special groups. In our investigations, the theory of multirings/hyperrings played a central role, and we got an interesting result: we have obtained an wide extension (to all special hyperfields, or special groups) of the validity of the Arason-Pfister Hauptsatz ([7]) - a positive answer to a question posed by J. Milnor in a classical paper of 1970 ([52])- and applied that to obtain interesting properties of graded rings associated to special hyperfields ([30, [18]).

Keywords: special groups, graded rings, profinite 2-group, multirings, hyperfield, quadratic forms, Arason-Pfister hauptsatz.

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## Introduction

It can be said that the Algebraic Theory of Quadratic Forms (ATQF) was founded in 1937 by E. Witt, with the introduction of the Witt ring concept of a given field, constructed from quadratic forms with coefficients in the field: given $F$, an arbitrary field of characteristic $\neq 2, W(F)$, the Witt Ring of $F$, classifies the quadratic forms on $F$ that are regular and anisotropic. Moreover, this ring establishes a strong connection between quadratic forms and orderings in a field $F$ : the set of orderings in $F$ is in one-to-one correspondence with the set of minimal prime ideals of the Witt ring of $F$, and more, the set of orderings in $F$ equipped with the Harrison's topology is a Boolean topological space and that by the bijection above, it is identified with a subspace of the Zariski spectrum of $W(F)$.

Further questions about Witt's ring structure $W(F)$ could only be answered about three decades after the original idea of Witt, through the introduction and analysis of the concept of Pfister form. The Pfister forms of degree $n \in \mathbb{N}$, in turn, are additive generators of the $n$-th power $I^{n}(F)$ of the fundamental ideal $I(F) \subseteq W(F)$ (the ideal determined by the anisotropic forms of even dimension).

In the early 1970s, J. Milnor, in his celebrated article [52], established deep functorial relationships, such as illustrated in the diagram below:


More specifically, Milnor determines a graded ring $k_{*}(F)$ (of reduced $K$-theory mod 2) associated to the field $F$, that interpolates, through morphisms of graded rings

$$
h_{*}(F): k_{*}(F) \longrightarrow H^{*}(F) \text { and } s_{*}(F): k_{*}(F) \longrightarrow W_{*}(F),
$$

Where

$$
\begin{gathered}
W_{*}(F):=\bigoplus_{n \in \mathbb{N}} I^{n}(F) / I^{n+1}(F) \\
H^{*}(F):=\bigoplus_{n \in \mathbb{N}} H^{n}\left(G a l\left(F^{s} \mid F\right),\{ \pm 1\}\right)
\end{gathered}
$$

are, respectively, the graded Witt ring of $F$ and the graded cohomology ring of $F$ (here, $F^{s}$ denotes the separable closure of $F$ ).

In this context, two fundamental questions are posed by Milnor:
i - Is true that $\bigcap_{n \in \mathbb{N}} I^{n}(F)=\{0\}$ ?
ii - Are the morphisms $h_{*}(F), s_{*}(F)$ isomorphisms of graded rings, for all field $F$ of characteristic not 2 ?

The question (i) was answered positively a few years later in a celebrated article by Arason and Pfister ([7]). Question (ii) resisted much longer until it was solved positively around the 2000s by V. Voevodsky and co-authors ([42]), results that earned to the first the Fields medal.

The absolute Galois group of $F, \operatorname{Gal}_{F}\left(F^{s}\right)$, detects the ordenability of $F$ : $F$ is formally real if and only if there is a non-trivial involution, i.e, an element $\sigma \in \operatorname{Gal}_{F}\left(F^{s}\right), \sigma \neq 1$ such that $\sigma^{2}=1$. But the Galois group detects more: the ordering space $X_{F}:=\operatorname{Sper}(F)$ is homeomorphic to the set

$$
\left\{[g]: g \text { is a non-trivial involution of } \operatorname{Gal}_{F}\left(F^{s}\right)\right\},
$$

where $[g]:=\left\{\sigma^{-1} g \sigma: \sigma \in \operatorname{Gal}_{F}\left(F^{s}\right)\right\}$.
We can see that, via the established functors, the Galois cohomology also describes orderings via the encoding of (graded) Witt rings. In fact, by Milnor's triangle, $W_{*}(F) \cong H^{*}\left(\operatorname{Gal}_{F}\left(F^{s}\right), \mathbb{Z}\right)$ and, keeping in mind that also $W_{*}(F)$ determines $W(F)$, we get a connection between the classic Witt ring $W(F)$ and Galois cohomology. Moreover, since the space of orderings $X_{F}$ is in natural bijection with $\operatorname{Hom}(W(F), \mathbb{Z})$, we obtain (again) a connection between orderings and Galois cohomology.

From the proof of Milnor's conjectures by Voevodski and the development of the theory of special groups - an abstract (first-order) theory of ATQF, introduced by M. Dickmann and F. Miraglia in the 1990s, - it was possible to demonstrate Conjectures on signatures put forward by M. Marshall and T. Lam in the mid-1970s ([27, [29]). We present below these two cases that exemplify the success of the application of instruments developed by the theory of the special groups:

Let $G$ be a special group, $X_{G}=\operatorname{Hom}_{S G}\left(G, \mathbb{Z}_{2}\right)$ the ordering space of $G$ and $W(G)$ the Witt ring associated to $G$. Denote $I^{n}(G)$ by the $n$-th power of the fundamental ideal (the ideal of forms of even dimension) in $W(G)$ and $W_{\text {tor }}(G)$ denote the torsion subgroup of $W(G)$. Consider the following statements:
a - For all $G$-form $\psi$, if $\operatorname{sgn}(\psi) \equiv 0 \bmod 2^{n}$ for all $\sigma \in X_{G}$, then $\psi \in I(G)$.
b - For all $G$-form $\psi$, if $\operatorname{sgn}_{\sigma}(\psi) \equiv 0 \bmod 2^{n}$ for all $\sigma \in X_{G}$, then $\psi \in I(G)+W_{\text {tor }}(G)$.
The Marshall's signature Conjecture ([MC]), originally stated in the (dual) context of abstract ordering spaces (see [48]), is that the statement (a) above is true for all reduced special group $G$. The Lam's Conjecture ([LC]), originally stated in the traditional context of formally real fields (see [43]), is that the statement (b) above is true for all formally real special group $G$. Note that both reverses of the above implications are true.

The fundamental stone for the solutions of Marshall's signature conjecture and Lam's conjecture for Pythagorean and formally real fields the introduction of the Boolean hull functor in SG-theory mainly through the definition of Horn-Tarski and Stiefel-Whitney invariants for isometry (see [28]). These encoding allows to rewrite $[\mathrm{MC}]$ in terms of information about the graded Witt ring and, through Voevodsky's results, switch these information in terms of the cohomology graded ring of a field, where the question is finally solved by the application of cohomological methods.

Entering in the "cohomology realm", in 1970's was established the first relations between quadratic forms and Galois cohomology, via the absolute Galois group $\operatorname{Gal}_{F}\left(F^{s}\right)([52])$. These
relations was improved, first via the quadratic closure $F^{q} \mid F$ in the 1980's, and in the late 1990's via a certain quotient

$$
\operatorname{Gal}_{F}\left(F^{q}\right) \rightarrow \operatorname{Gal}_{F}\left(F^{(3)}\right)
$$

called the $W$-group of $F$ (see [53]). Moreover, the induced arrow

$$
H^{*}\left(\operatorname{Gal}\left(F^{s} \mid F\right),\{ \pm 1\}\right) \rightarrow H^{*}\left(\operatorname{Gal}\left(F^{(3)} \mid F\right),\{ \pm 1\}\right)
$$

is a monomorphism whose image is the subgraded ring of $H^{*}\left(\operatorname{Gal}\left(F^{(3)} \mid F\right),\{ \pm 1\}\right)$ generated by cup products of level 1 members. However, only after 2010 the studies of this setting ${ }^{1}$ have established that $F^{(3)} \mid F$ is the minimal extension that determines (and is determined by) $W(F)$.

All that was exposed above compose an amount of evidences so that we propose a new diagram


Our proposal is to investigate the precise mathematical relationships so that the above diagram will be true, possible configuring a "new adjoint" situation for these theories, obviously with the intention of exploring the possible transport of information from this. It is important to point out that the current paradigm in abstract quadratic forms theories is that of equivalence (or duality) of categories, which is a relatively rigid connection. The context of adjunction allows more flexibility in attacking problems from one theory encoded in another.

In Chapter 1 we present the theory of multirings and multifields/hyperfields, closely to the perspective of Marshall's paper [47]. Roughly speaking, multirings are just "rings with a multivalued addition". In fact, many ideas of the ring theory can be imported. The main references are [47, [24], [23], [45] and [58, and we follow 47] closely. In fact, the main proofs concerning orderings over hyperfields are easier than the field case, and we got an Artin-Schreier Theory very similar to the field one (see for instance, Propositions 1.4 .5 and 1.6.3). We also characterize real reduced hyperfields (Corollary 1.5.3) and real reduced multirings (Proposition 1.8.4), which are respectively dual to Marshall's abstract ordering spaces and abstract real spectra.

In Chapter 2 we connect the theory of Chapter 1 with the with the most significant theories of quadratic forms, via two main motivations: 1 ) to describe interesting pairs $(A, T)$ where $A$ is a (multi)ring and $T \subseteq A$ is a certain multiplicative subset in such a way to obtain models of abstract theories of quadratic forms (special groups and real semigroups) via natural quotients - Marshall's quotient construction; and 2) use this construction to motivate a "non reduced" expansion of the theory of real semigroups to deal the formally real case, isolating axioms over pairs involving multirings and a subset with some properties. The main results are Theorem 2.3.4, 2.3.7, 2.3.10 and 2.5.4 which characterize precisely the necessary conditions for a hyperfield/multiring come from a

[^0]special group/real semigroup. Proposition 2.4.3 deals with a question posed by the authors of 37. We also got a new and interesting Example of real semigroup 2.5.15): $A /{ }_{m} T$ for $A=\mathcal{C}(X, \mathbb{R})$ and $T=A^{2} \cap \operatorname{nzd}(A)$, where $X$ is a $T_{6}$ topological space.

In Chapter 3 we introduce the theory of superrings. They are important in order to obtain the quadratic extension available for special groups. The concept of superring first appears in ([6]). There are many important advances and results in hyperring theory, and for instance, we recommend for example, the following papers: [3], [5, [6, [4], 49], [54], 51, [50]. Surprisingly we have obtained an interesting theory of matrices, linear systems, vector spaces and algebraic extensions available for a certain subclass of superfields. If $R$ is a full superring, then $M_{m \times n}(R)$ and $R[X]$ are superrings (Theorem 3.2.6 and 3.4.2). We also obtained a kind of simple algebraic extension for a superfield $F$ (Theorem 3.6.12), which culminate in the existence and unicity of a full algebraic extension of a superfield $F$ (Theorems 3.7 .3 and 3.7.4. If $F$ is a linearly closed superfield (the system $A x=0$ always have a non trivial solution), then we have a well defined dimension theory for the vector spaces over $F$ (Theorem 3.8.21). The main examples of linearly closed superfields are hyperbolic hyperfields (3.8.23) and simple full algebraic extensions over a linearly closed superfield (3.8.25). The linearly closed interpreted in the context of special groups leads to interesting Isotropic (Corollary 3.8.27) and Hyperbolic (Corollary 3.8.28) interpolations. We finish this Chapter with a quantifier elimination procedure for superfields (Theorem 3.9.3), which is a direct generalization of a result obtained in [19.

In Chapter 4 we we provide some new steps towards the development of tools of algebraic theory of quadratic forms in this multiring setting: we have defined and explored K-theory and graded rings in the context of hyperfields that, in particular, provides a generalization and unification of Milnor's K-theory ([52]) and special groups K-theory ( 30$]$ ). We develop some properties of this generalized K-theory, that can be seen as a free inductive graded ring. The main results are Theorems 4.5 .6 and its Corollaries, which provides interchanging formulas between the three K-theories considered here.

In Chapter 5 we deal with the category IGR. Theorem 4.5 .6 gives a hint that the category of Igr is a good abstract environment for studying questions of "quadratic flavour". So a better understanding of Igr's is at least desirable and this is the main purpose of this Chapter. We develop the general properties valid for Igr's and the main results here are Theorem 5.5.4, providing an adjunction between the categories of pre-special groups and (a subcategory of) inductive graded rings. We also characterize the Special and Weak Marshall Conjecture in the context of inductive graded rings (Section 5.6).

In Chapter 6 we develop the theory of quadratic extensions for hyperfields/superfields, through the development of results concerning the superrings of polynomials, envisaging some applications to algebraic theory of quadratic forms and Real Algebraic Geometry. The main results here are the Arason-Pfister Hauptsatz for all special groups (Theorem 6.3.2) and its consequences.

The Igr's functors $W_{*}, k_{*}$ were extended by M. Dickmann and F.Miraglia from the category of fields of characteristic $\neq 2$ to the category of special groups (equivalently, the category of special hyperfields). Another relevant Igr functor, the graded cohomology ring, $H^{*}\left(\operatorname{Gal}\left(F^{s} \mid F\right),\{ \pm 1\}\right)$ remains defined only on the field setting. Chapter 7 constitutes an attempt to provide an Igr functor associated to a (Galois) cohomology theory for special groups, based on the work of J. Minac and M. Spira [53]: we will define - by "generator and relations", $\operatorname{Gal}(G)$, the Galois Group of an $S G G$, and provide some properties of this construction, as the encoding of the orderings on $G$. However, since deeper results will depend of a description of $\operatorname{Gal}(G)$ "from below", and it still unavailable a complete theory of algebraic extension of (super)hyperfields, we will not pursue a more complete development of this cohomology theory in this thesis, reserving it for a future research. The main results are Theorem 7.3.13 and 7.3.15, which recover for the abstract context
the characterization of orderings in terms of the involutions in the Galois group of a field.
In Chapter 8 we finish the work indicating some possibilities of future research connected with this thesis.

## Chapter 1

## Multirings and Hyperfields

Here, we present the theory of multirings and multifields/hyperfields, closely to the perspective of Marshall's paper [47]. Roughly speaking, multirings are just "rings with a multivalued addition". In fact, many ideas of the ring theory can be imported. The main references are [47, [24, [23], [45] and [58], and we follow [47] closely.

In fact, the main proofs concerning orderings over hyperfields are easier than the field case, and we got an Artin-Schreier Theory very similar to the field one (see for instance, Propositions 1.4.5 and 1.6.3.

We also characterize real reduced hyperfields (Corollary 1.5.3) and real reduced multirings (Proposition 1.8.4), which are respectively dual to Marshall's abstract ordering spaces and abstract real spectra.

### 1.1 On Multialgebras

There are several Definitions of multialgebra on the literature, considering that each multialgebra application in a specific area of Mathematics (mainly Algebra and Logic) requires a particular adaptation. Here, we adapt the notion of multialgebra used in [10]; the identity theory here presented is close to the exposed in 55].

Definition 1.1.1. A multialgebraic signature is a sequence of pairwise disjoint sets

$$
\Sigma=\left(\Sigma_{n}\right)_{n \in \mathbb{N}}
$$

where $\Sigma_{n}=S_{n} \sqcup M_{n}$, which $S_{n}$ is the set of strict multi-operation symbols and $M_{n}$ is the set of multioperation symbols. In particular, $\Sigma_{0}=S_{0} \sqcup M_{0}, F_{0}$ is the set of symbols for constants and $M_{0}$ is the set of symbols for multi-constants. We also denote

$$
\Sigma=\left(\left(S_{n}\right)_{n \geq 0},\left(M_{n}\right)_{n \geq 0}\right)
$$

Definition 1.1.2. Let $A$ be any set.
$i$ - A multi-operation of arity $n \in \mathbb{N}$ over a set $A$ is a function

$$
A^{n} \rightarrow \mathcal{P}^{*}(A):=\mathcal{P}(A) \backslash\{\emptyset\} .
$$

ii - A multi-operation of arity $n \in \mathbb{N}$ over a set $A, A^{n} \rightarrow \mathcal{P}^{*}(A)$, is strict, whenever it factors through the singleton function $s_{A}: A \mapsto \mathcal{P}^{*}(A), a \mapsto s_{A}(a):=\{a\}$. Thus it can be naturally identified with an ordinary n-ary operation $A^{n} \rightarrow A$.

A 0 -ary multi-operation (respectively strict multi-operation) on $A$ can be identified with a non-empty subset of $A$ (respectively a singleton subset of $A$ ).

Definition 1.1.3. A multialgebra over a signature $\Sigma=\left(\left(S_{n}\right)_{n \geq 0},\left(M_{n}\right)_{n \geq 0}\right)$, is a set $A$ endowed with a family of n-ary multioperations

$$
\sigma_{n}^{A}: A^{n} \rightarrow \mathcal{P}^{*}(A), \sigma_{n} \in S_{n} \sqcup M_{n}, n \in \mathbb{N},
$$

such that: if $\sigma_{n} \in S_{n}$, then $\sigma_{n}^{A}: A^{n} \rightarrow \mathcal{P}^{*}(A)$ is a strict $n$-ary multioperation.

## Remark 1.1.4.

$i$ - Every algebraic signature $\Sigma=\left(F_{n}\right)_{n \in \mathbb{N}}$ is a multialgebraic signature where $M_{n}=\emptyset$, for all $n \in \mathbb{N}$. Each algebra

$$
\left(A,\left(\left(A^{n} \xrightarrow{f^{A}} A\right)_{f \in F_{n}}\right)_{n \in \mathbb{N}}\right)
$$

over the algebraic signature $\Sigma$ can be naturally identified with a multi-algebra

$$
\left(A,\left(\left(A^{n} \xrightarrow{f^{A}} A \xrightarrow{s_{A}} \mathcal{P}^{*}(A)\right)_{f \in F_{n}}\right)_{n \in \mathbb{N}}\right)
$$

over the same signature.
ii - Every multialgebraic signature $\Sigma=\left(\left(S_{n}\right)_{n \in \mathbb{N}},\left(M_{n}\right)_{n \in \mathbb{N}}\right)$ induces naturally a first-order language

$$
L(\Sigma)=\left(\left(F_{n}\right)_{n \in \mathbb{N}},\left(R_{n+1}\right)_{n \in \mathbb{N}}\right)
$$

where $F_{n}:=S_{n}$ is the set of $n$-ary operation symbols and $R_{n+1}:=M_{n}$ is the set of $(n+1)$-ary relation symbols. In this way, multi-algebras

$$
\left(A,\left(\left(A^{n} \xrightarrow{\sigma^{A}} \mathcal{P}^{*}(A)\right)_{\sigma \in S_{n} \sqcup M_{n}}\right)_{n \in \mathbb{N}}\right)
$$

over a multialgebraic signature $\Sigma=\left(S_{n} \sqcup M_{n}\right)_{n \in \mathbb{N}}$ can be naturally identified with the firstorder structures over the language $L(\Sigma)$ that satisfies the $L(\Sigma)$-sentences:

$$
\forall x_{0} \cdots \forall x_{n-1} \exists x_{n}\left(\sigma_{n}\left(x_{0}, \cdots, x_{n-1}, x_{n}\right)\right), \text { for each } \sigma_{n} \in R_{n+1}=M_{n}, n \in \mathbb{N} .
$$

Now we focus our attention into a more syntactic aspect of this multi-algebras theory. We start with a (recursive) definition of multi-terms:

Definition 1.1.5. $A$ (multi-)term on a multialgebra $A$ of signature

$$
\Sigma=\left(\left(S_{n}\right)_{n \geq 0},\left(M_{n}\right)_{n \geq 0}\right)
$$

is defined recursively as:
$i$ - Variables $x_{i}, i \in \mathbb{N}$ are terms.
ii - If $t_{0}, \cdots, t_{n-1}$ are terms and $\sigma \in S_{n} \sqcup M_{n}$, then $\sigma\left(t_{0}, \cdots, t_{n-1}\right)$ is a term.
We will call a multi-term $t$ strict, whenever it is composed only by combination of strict multioperations and variables. The notion of occurrence of a variable in a term is as the usual. We will denote var $(t)$ as the (finite set of variables) that occurs in the term $t$.

To define an interpretation for terms, we need a preliminary step. Given

$$
\sigma \in S_{n} \sqcup M_{n},
$$

we "extend" $\sigma^{A}: A^{n} \rightarrow \mathcal{P}^{*}(A)$ to a n-ary operation in $\mathcal{P}^{*}(A)$,

$$
\sigma^{\mathcal{P}^{*}(A)}: \mathcal{P}^{*}(A)^{n} \rightarrow \mathcal{P}^{*}(A)
$$

by the rule:

$$
\sigma^{\mathcal{P}^{*}(A)}\left(A_{0}, \cdots, A_{n-1}\right):=\bigcup_{a_{0} \in A_{0}} \cdots \bigcup_{a_{n-1} \in A_{n-1}} \sigma^{A}\left(a_{0}, \cdots, a_{n-1}\right) .
$$

Definition 1.1.6. The interpretation of a term $t$ on a multialgebra $A$ over a signature

$$
\Sigma=\left(\left(S_{n}\right)_{n \geq 0},\left(M_{n}\right)_{n \geq 0}\right)
$$

is a function $t^{A}: A^{\operatorname{var}(t)} \rightarrow \mathcal{P}^{*}(A)$ and is defined recursively as follows:
$i$ - The interpretation of a variable $x_{i}, x_{i}^{A}: A^{\left\{x_{i}\right\}} \rightarrow \mathcal{P}^{*}(A)$ is essentially the singleton function of $A$ :

$$
x_{i}^{A}: A^{\left\{x_{i}\right\}} \cong A \mapsto \mathcal{P}^{*}(A) \text {, is given by the rule }\left(\hat{a}:\left\{x_{i}\right\} \rightarrow A\right) \mapsto\{a\} .
$$

ii - If $t=\sigma\left(t_{0}, \cdots, t_{n-1}\right)$ is a term and $\sigma \in S_{n} \sqcup M_{n}$, denote $T=\operatorname{var}(t)$ and $T_{i}=\operatorname{var}\left(t_{i}\right)$. Then $T=\bigcup_{i<n} T_{i}$. Consider $t_{i}^{A}: A^{T} \rightarrow \mathcal{P}^{*}(A)$ the composition

$$
A^{T} \xrightarrow{\operatorname{proj}_{T_{i}}^{T}} A^{T_{i}} \xrightarrow{t^{A}} \mathcal{P}^{*}(A)
$$

where $\operatorname{proj}_{T_{i}}^{T}$ is the canonical projection induced by the inclusion $T_{i} \hookrightarrow T$. Then

$$
t^{A}: A^{T} \rightarrow \mathcal{P}^{*}(A)
$$

is the composition

$$
A^{T} \xrightarrow{\left.\left(t_{i}\right)_{T}\right)_{i<n}}\left(\mathcal{P}^{*}(A)\right)^{n} \xrightarrow{\mathcal{P}^{*}(A)} \mathcal{P}^{*}(A) .
$$

Definition 1.1.7. Let $A$ be a multialgebra $A$ over a signature $\Sigma=\left(\left(S_{n}\right)_{n \geq 0},\left(M_{n}\right)_{n \geq 0}\right)$ and let $t_{1}, t_{2}$ be $\Sigma$-terms. We say that $A$ realize that $t_{1}$ is contained in $t_{2}$, (notation: $A \models t_{1} \sqsubseteq t_{2}$ ) whenever $t_{1}^{A}(\bar{a}) \subseteq t_{2}^{A}(\bar{a})$, for each tuple $\bar{a}: \operatorname{var}\left(t_{1}\right) \cup \operatorname{var}\left(t_{2}\right) \rightarrow A$.

Apart from the notion of atomic formulas the definition of $\Sigma$-formulas for multi-algebraic theories is similar to the (recursive) definition of first-order $L(\Sigma)$-formulas:
Definition 1.1.8. The formulas of $\Sigma$ are defined as follows:
$i$ - Atomic formulas are the formulas of type $t \sqsubseteq t^{\prime}$, where $t, t^{\prime}$ are terms.
ii- If $\phi, \psi$ are formulas, then $\neg \phi$ and $\phi \vee \psi, \phi \wedge \psi, \phi \rightarrow \psi, \phi \leftrightarrow \psi$ are formulas.
iii- If $\phi$ is a formula and $x_{i}$ is a variable, then $\forall x_{i} \phi, \exists x_{i} \phi$ are formulas
The notion of occurrence (respectively free occurrence) of a variable in a formula is as the usual. We will denote $f v(\phi)$ as the (finite) set of variables that occurs free in the formula $\phi$.

We use $t_{1}={ }_{s} t_{2}$ to abbreviate the formula $\left(t_{1} \sqsubseteq t_{2}\right) \wedge\left(t_{2} \sqsubseteq t_{1}\right)$ : this means that $t_{1}$ and $t_{2}$ are "strongly equal terms".

Definition 1.1.9. The definition of interpretation of formulas $\phi(\bar{x})$ where

$$
f v(\phi) \subseteq \bar{x} \subseteq\left\{x_{i}: i \in \mathbb{N}\right\}
$$

under a valuation of variables $v: \bar{x} \rightarrow A$ (or we will denote simply by $v=\bar{a}$ ) is:

$$
i-A \models_{v} t(\bar{x}) \sqsubseteq t^{\prime}(\bar{x}) \text { iff } t^{A}(\bar{a}) \subseteq t^{\prime A}(\bar{a})
$$

ii- The case of complex formulas (given by the connectives $\neg, \vee, \wedge, \rightarrow, \leftrightarrow$, and quantifiers $\forall, \exists$ ) is as satisfaction of first-order $L(\Sigma)$-formulas in $L(\Sigma)$-structure on a valuation $v$.

## Remark 1.1.10.

$i$ - The theory of multi-algebras entails that for each term $t$, and each strict term $t^{\prime}$,

$$
t \sqsubseteq t^{\prime} \text { iff } t={ }_{s} t^{\prime}
$$

ii- In [55] contains a development of the identity theory for multialgebras, with another primitive notion: $t(\bar{x})={ }_{w} t^{\prime}(\bar{x})$; a $\Sigma$-multialgebra A satisfies the "weak identity" above iff there is some $\bar{a} \in A^{\operatorname{var}(t) \cup v a r\left(t^{\prime}\right)}$ such that $t^{A}(\bar{a}) \cap t^{\prime A}(\bar{a}) \neq \emptyset$. This will not play any role in this work but is useful for applications of multi-algebraic semantics for complex logical systems ([38]).

There are many ways of define morphism for multialgebras. Follow below our choice:
Definition 1.1.11. Let $A$ and $B$ be multialgebras of signature $\Sigma=\left(\left(S_{n}\right)_{n \geq 0},\left(M_{n}\right)_{n \geq 0}\right)$ and $\varphi: A \rightarrow B$ be a function.
$i-\varphi$ is a partial morphism if for every $n \geq 0$, every $\sigma \in S_{n}$ and every $a_{1}, \ldots, a_{n} \in A$, we have

$$
\varphi\left(\sigma^{A}\left(a_{1}, \ldots, a_{n}\right)\right) \subseteq \sigma^{B}\left(\varphi\left(a_{1}\right), \ldots, \varphi\left(a_{n}\right)\right)
$$

ii - $\varphi$ is a morphism if for every $n \geq 0$, every $\sigma \in S_{n} \sqcup M_{n}$ and every $a_{1}, \ldots, a_{n} \in A$, we have

$$
\varphi\left(\sigma^{A}\left(a_{1}, \ldots, a_{n}\right)\right) \subseteq \sigma^{B}\left(\varphi\left(a_{1}\right), \ldots, \varphi\left(a_{n}\right)\right)
$$

iii - $\varphi$ is a strong morphism if for every $n \geq 0$, every $\sigma \in S_{n} \sqcup M_{n}$ and every $a_{1}, \ldots, a_{n} \in A$, we have

$$
\varphi\left(\sigma^{A}\left(a_{1}, \ldots, a_{n}\right)\right)=\sigma^{B}\left(\varphi\left(a_{1}\right), \ldots, \varphi\left(a_{n}\right)\right)
$$

## Remark 1.1.12.

$i$ - Let $A, B$ be $\Sigma$-multialgebras. If $B$ is a strict multilagebra (i.e. $\sigma_{n}^{B}(\bar{b})$ is unitary subset of $B$, for each $\sigma \in \Sigma$ and each tuple $\bar{b}$ in $B)$, then the morphisms $A \rightarrow B$ coincide with the strong morphisms $A \rightarrow B$.
ii - There is a full and faithful concrete embedding of the category of ordinary algebraic structures over a signature $\Sigma$ and homomorphisms into the category of $\Sigma$-multialgebras and (strong) morphisms: the image of this embedding is the class of strict multialgebras over $\Sigma$.
iii - The correspondence $\Sigma \mapsto L(\Sigma)$ induces a concrete isomorphism between the category of $\Sigma$ multialgebras and the category of $L(\Sigma)$ - first order structures satisfying suitable $\forall \exists$ axioms. It is ease to see that this correspondence induces a bijection between injective strong embedding of $\Sigma$-multialgebras and $L(\Sigma)$-monomorphisms of first-order structures.

We finish this subsection with two illustrative examples of multialgebras derived from an algebraic structure and from a first-order structure.

Example 1.1.13. Let $(R,+, \cdot, 0,1)$ be a commutative ring with $1 \neq 0$. Given $n \geq 1$, define an $(n+1)$-ary multioperation $*_{n}$ by the rule:

$$
\begin{aligned}
d \in a_{0} *_{n} a_{1} *_{n} a_{2} *_{n} \ldots *_{n} a_{n} \Leftrightarrow \text { there is some } t & \in R \text { such that } \\
d & =a_{0}+a_{1} t+a_{2} t^{2}+\ldots+a_{n} t^{n} .
\end{aligned}
$$

The idea here, is that $a_{0} *_{n} a_{1} *_{n} a_{2} *_{n} \ldots *_{n} a_{n}$ "analyze" the values taken in $R$ by the polynomial $p(X)=a_{0}+a_{1} X+a_{2} X^{2}+\ldots+a_{n} X^{n} \in R[X] . *_{n}$ will be called The streching multialgebra of degree $n$ over $R$.
ordermulti
Example 1.1.14. Let $\mathcal{L}=\{0,1,+, \cdot, \leq\}$ the language of ordered fields. Consider $\mathbb{R}$ as an ordered field. We can look at the ordering relation as a multioperation of arity 1. In agreement with our notation, we have

$$
\leq(a):=\{x \in \mathbb{R}: a \leq x\}=[a,+\infty) .
$$

From now on, all multi-algebras considered in this work will contain only operations of arities $0,1,2$. They will have strict constants and strict unary operations; the binary operations maybe strict or multivalued.

### 1.2 Multigroups, Multirings and Multifields

Multigroups are a generalization of groups. We can think that a multigroup is a group with a multivalued operation:
defn:multigroupI
Definition 1.2.1. A multigroup is a quadruple $(G, *, r, 1)$, where $G$ is a non-empty set, functions $*: G \times G \rightarrow \mathcal{P}(G) \backslash\{\emptyset\}$ and $r: G \rightarrow G$ and 1 is an element of $G$ satisfying:
$i$ - If $z \in x * y$ then $x \in z * r(y)$ and $y \in r(x) * z$.
$i i-y \in 1 * x$ iff $x=y$.
iii - With the convention $x *(y * z)=\bigcup_{w \in y * z} x * w$ and $(x * y) * z=\bigcup_{t \in x * y} t * z$,

$$
x *(y * z)=(x * y) * z \text { for all } x, y, z \in G .
$$

A multigroup is said to be commutative if
iv $-x * y=y * x$ for all $x, y \in G$. The structure $(G, \cdot, 1)$ is a commutative multimonoid (with unity) if satisfy M3 and M4 and the condition $a \in 1 \cdot a$ for all $a \in G$.
mgmorph1
Definition 1.2.2. A morphism of multigroups is a function $f: G \rightarrow H$ between multigroups such that $f\left(1_{G}\right)=1_{H}$ and for all $a, b \in G, f(a * b) \subseteq f(a) * f(b)$. Of course, composition of morphisms is a morphism and the category of multigroups with their morphisms will be denoted by MGrp.

There is another definition (due to Marshall in 47]) with a first order theoretic flavour.

Definition 1.2.3 (Adapted from Definition 1.1 in [47]). A multigroup is a quadruple ( $G, \Pi, r, \mathfrak{i}$ ) where $G$ is a non-empty set, $\Pi$ is a subset of $G \times G \times G, r: G \rightarrow G$ is a function and $\mathfrak{i}$ is an element of $G$ satisfying:

I-If $(x, y, z) \in \Pi$ then $(z, r(y), x) \in \Pi$ and $(r(x), z, y) \in \Pi$.
$I I-(x, \mathfrak{i}, y) \in \Pi$ iff $x=y$.
III - If there exist $p \in G$ such that $(u, v, p) \in \Pi$ and $(p, w, x) \in \Pi$ then there exist $q \in G$ such that $(v, w, q) \in \Pi$ and $(u, q, x) \in \Pi$.
A multigroup is said to be commutative if
$I V-(x, y, z) \in \Pi$ iff $(y, x, z) \in \Pi$.
In fact, these Definitions describes the same object, and that connection is established by the following Lemma:

Lemma:1.2
Lemma 1.2.4 (Lemma 1.3 of [47]). For any multigroup $G$ as in the second version 1.2 .3 , we have:
$a-r(\mathfrak{i})=\mathfrak{i}$.
$b-r(r(x))=x$.
$c-(x, y, z) \in \Pi i f f(r(y), r(x), r(z)) \in \Pi$.
$d-(\mathfrak{i}, x, y) \in \Pi$ iff $x=y$.
$e$ - If there exist $q \in G$ such that $(v, w, q) \in \Pi$ and $(u, q, x) \in \Pi$ then there exist $p \in G$ such that $(u, v, p) \in \Pi$ and $(p, w, x) \in \Pi$.
$f$ - For each $a, b \in G$, there exists $c \in G$ such that $(a, b, c) \in \Pi$.
Proof.
a - Since $i=i$, by II we have $(i, i, i) \in \Pi$. By I, $(r(i), i, i) \in \Pi$ and by II, $r(i)=i$.
$\mathrm{b}-x=x \stackrel{I I}{\Rightarrow}(x, \mathfrak{i}, x) \in \Pi \stackrel{I}{\Rightarrow}(r(x), x, \mathfrak{i}) \in \Pi \stackrel{I}{\Rightarrow}(r(r(x)), \mathfrak{i}, x) \in \Pi \stackrel{I I}{\Rightarrow} r(r(x))=x$.
c - $(x, y, z) \stackrel{I}{\Leftrightarrow}(z, r(y), x) \in \Pi \stackrel{I}{\Leftrightarrow}(r(z), x, r(y)) \in \Pi \stackrel{I}{\Leftrightarrow}(r(y), r(x), r(z)) \in \Pi$.
d - Let $(\mathfrak{i}, x, y) \in \Pi$. Then

$$
\begin{array}{r}
(\mathfrak{i}, x, y) \in \Pi \stackrel{I}{\Rightarrow}(y, r(x), \mathfrak{i}) \in \Pi \stackrel{I}{\Rightarrow}(r(y), \mathfrak{i}, r(x)) \in \Pi \\
\quad \Rightarrow r(y)=r(x) \stackrel{(b)}{\Rightarrow} y=r(r(y))=r(r(x))=x .
\end{array}
$$

Conversely, suppose $x=y$. Then

$$
\begin{aligned}
x=y & \Rightarrow r(x)=r(y) \\
& \stackrel{I I}{\Rightarrow}(r(y), \mathfrak{i}, r(x)) \in \Pi \\
& I(b) \\
\Rightarrow & (y, r(x), \mathfrak{i}) \in \Pi \stackrel{I}{\Rightarrow}(\mathfrak{i}, x, y) \in \Pi .
\end{aligned}
$$

e- Note that

$$
(u, q, x) \in \Pi \stackrel{I}{\Rightarrow}(x, r(q), u) \in \Pi \stackrel{(c)}{\Rightarrow}(q, r(x), r(u)) \in \Pi .
$$

Then, $(v, w, q) \in \Pi$ and $(q, r(x), r(u)) \in \Pi$, so by axiom III, there exists $t \in G$ such that $(w, r(x), t) \in \Pi$ and $(v, t, r(u)) \in \Pi$.

$$
\begin{aligned}
(w, r(x), t) & \in \Pi \stackrel{(b)}{\Rightarrow}(x, r(w), t) \in \Pi \stackrel{I}{\Rightarrow}(r(t), w, x) \in \Pi, \text { and } \\
(v, t, r(u)) & \in \Pi \stackrel{(b)}{\Rightarrow}(r(t), r(v), u) \in \Pi \stackrel{I}{\Rightarrow}(u, v, r(t) \in \Pi .
\end{aligned}
$$

Hence defining $p=r(t)$, we have $(u, v, p) \in \Pi$ and $(p, w, x) \in \Pi$.
f - Since $(b, r(b), \mathfrak{i}) \in \Pi$ and $(a, \mathfrak{i}, a) \in \Pi$, by (e), there exists $c \in G$ such that $(a, b, c) \in \Pi$ and $(c, r(b), a) \in \Pi$.
mgmorph2
Definition 1.2.5. A morphism of multigroups (in the sense of Definition 1.2.3) is a function $f: G \rightarrow H$ between multigroups $\left(G, \Pi_{G}, r_{G}, \mathfrak{i}_{G}\right)$ and $\left(H, \Pi_{H}, r_{H}, \mathfrak{i}_{H}\right)$ such that $f\left(1_{G}\right)=f\left(1_{H}\right)$ and for all $a, b, c \in G$ if $(a, b, c) \in \Pi_{G}$ then $(f(a), f(b), f(c)) \in \Pi_{H}$. Of course, composition of morphisms is a morphism and the category of multigroups with their morphisms will be denoted by $M G r p_{f o l}$. ${ }^{\text {1 }}$

Theorem 1.2.6. The categories $M G r p$ and $M G r p_{f o l}$ are equivalent.
Proof. Let $(G, *, r, 1)$ be an object of $M G r p$. Define a multigroup $G_{f o l}:=\left(G, \Pi_{*}, r, \mathfrak{i}\right)$ taking $\mathfrak{i}=1$ and $\Pi_{*}=\{(a, b, c): c \in a * b\}$. The validity of axioms I,II, III (and IV) for $G_{f o l}$ are direct consequence of axioms i,ii, iii (and iv) for ( $G, *, r, 1$ ), so $G_{f o l}$ is an object in MGrp fol .

Conversely, let ( $G, \Pi, r, \mathfrak{i}$ ) be an object in $M G r p_{f o l}$. By $1.2 .4(\mathrm{f})$, we have a well-defined function $*_{\Pi}: A \times A \rightarrow \mathcal{P}(A) \backslash\{\emptyset\}$, given by the rule

$$
*_{\Pi}(a, b)=a *_{\Pi} b:=\{c \in G:(a, b, c) \in \Pi\} .
$$

Let $G_{M}:=\left(G, *_{\Pi}, 1\right)$ with $1=\mathfrak{i}$. Then, the validate of the axioms i,ii (and iv) for $G_{M}$ are direct consequence of I,II (and IV) for ( $G, \Pi, r, \mathfrak{i}$ ). For the axiom iii, let $x \in a *_{\Pi}\left(b *_{\Pi} c\right)$. Then $x \in a *_{\Pi} q$ for some $q \in b *_{\Pi} c$. Since $(b, c, q) \in \Pi$ and $(a, q, x) \in \Pi$, by $1.2 .4(\mathrm{e})$, there exists $p \in \Pi$ such that $(a, b, p) \in \Pi$ and $(p, c, x) \in \Pi$ and then, $x \in p *_{\Pi} c$ with $p \in a *_{\Pi} b$ that imply $x \in\left(a *_{\Pi} b\right) *_{\Pi} c$. Finally, let $y \in\left(a *_{\Pi} b\right) *_{\Pi} c$. So $y \in p *_{\Pi} c$ for some $p \in a *_{\Pi} b$, then and $(a, b, p) \in \Pi$ and $(p, c, y) \in \Pi$. By III, there exists $q \in \Pi$ such that $(b, c, q) \in \Pi$ and $(a, q, y) \in \Pi$. Hence $y \in a *_{\Pi} q$ and $q \in b *_{\Pi} c$, that imply $y \in a *_{\Pi}\left(b *_{\Pi} c\right)$. Therefore, $G_{M}$ is an object in MGrp.

Using the above arguments, we have the equivalence of these categories witnessed by the functors $\mathcal{F}: M G r p \rightarrow M G r p_{f o l}$ and $\mathcal{G}: M G r p_{f o l} \rightarrow M G r p$ defined respectively on the objects by $\mathcal{F}(G)=G_{f o l}$ and $\mathcal{G}(G)=G_{M}$; and on the morphisms $f \in \operatorname{MGrp}(G, H)$ and $g: \operatorname{MGrp}_{f o l}(K, L)$ by $\mathcal{F}(f)=f$ and $\mathcal{G}(g)=g$.

Now we deal with multirings.
Definition 1.2.7 (Adapted from Definition 2.1 in [47]). A multiring is a sextuple ( $R,+, \cdot,-, 0,1$ ) where $R$ is a non-empty set, $+: R \times R \rightarrow \mathcal{P}(R) \backslash\{\emptyset\}, \cdot: R \times R \rightarrow R$ and $-: R \rightarrow R$ are functions, 0 and 1 are elements of $R$ satisfying:

[^1]$i-(R,+,-, 0)$ is a commutative multigroup;
ii - $(R, \cdot, 1)$ is a monoid;
iii $-a 0=0$ for all $a \in R$;
$i v$ - If $c \in a+b$, then $c d \in a d+b d$ and $d c \in d a+d b$. Or equivalently, $(a+b) d \subseteq a b+b d$ and $d(a+b) \subseteq d a+d b$.
$v$ - If the equalities holds, i.e, $(a+b) d=a b+b d$ and $d(a+b)=d a+d b$, we said that $R$ is $a$ hyperring.

A multiring is commutative if $(R, \cdot, 1)$ is a commutative monoid. A zero-divisor of a multiring $R$ is a non-zero element $a \in R$ such that $a b=0$ for another non-zero element $b \in R$. The multiring $R$ is said to be a multidomain if do not have zero divisors, and $R$ will be a multifield if $1 \neq 0$ and every non-zero element of $R$ has multiplicative inverse.

Remark 1.2.8. It is straightforward to realize that every multifield $F$ is in fact a hyperfield, i.e, for all $a, b, d \in F, d(a+b)=d a+d b$.

> ex:1.3

## Example 1.2.9.

$a$ - Suppose $(G, \cdot, 1)$ is a group. Defining $*(a, b)=\{c \in G: c=a \cdot b\}$ and $r(g)=g^{-1}$, we have that ( $G, *, r, 1$ ) is a multigroup.
$b$ - In the same way of item (a), every ring, domain and field is a multiring, multidomain and multifield respectively.
$\left.c-Q_{2}=\{-1,0,1\}\right]^{2}$ is a multifield with the usual product and the multivalued sum defined by relations

$$
\left\{\begin{array}{l}
0+x=x+0=x, \text { for every } x \in Q_{2} \\
1+1=1,(-1)+(-1)=-1 \\
1+(-1)=(-1)+1=\{-1,0,1\}
\end{array}\right.
$$

$d$ - Let $K=\{0,1\}$ with the usual product and the sum defined by relations $x+0=0+x=x$, $x \in K$ and $1+1=\{0,1\}$. This is a multifield called Krasner's multifield [41].

Example 1.2.10 (Example 2.5 of [47). Let be $V \subseteq \mathbb{R}^{n}$ an algebraic set and $A$ as the coordinate ring of $V$, i.e, the ring $\mathbb{R}[V]$ of polynomial functions $f: V \rightarrow \mathbb{R}$. Define an equivalence relation $\sim$ on $A$ by $f \sim g$ iff $f(x)$ and $g(x)$ has the same sign for all $x \in V$. Thus, $Q_{r e d}(A)=A / \sim$ is called the real reduced multiring. The operations are defined by:

$$
\left\{\begin{array}{l}
\bar{f} \in \bar{g}+\bar{h} \Leftrightarrow \exists f^{\prime}, g^{\prime}, h^{\prime} \in A \\
\quad \text { such that } f^{\prime}=g^{\prime}+h^{\prime}, \overline{f^{\prime}}=\bar{f}, \overline{g^{\prime}}=\bar{g}, \text { and } \overline{h^{\prime}}=\bar{h} \\
\bar{g} \bar{h}=\overline{g h},-\bar{f}=\overline{-f}, 0=\overline{0}, 1=\overline{1}
\end{array}\right.
$$

Taking $n=1$, we have a counter-example to show that $a d+b d \subsetneq(a+b) d$ in general:

$$
\overline{x^{2}+x^{3}} \in \overline{x x}+\bar{x} \overline{1} \text { but } \overline{x^{2}+x^{3}} \notin \bar{x}(\bar{x}+\overline{1}),
$$

and this not happen because $x^{2}+x^{3}>0$ and $x(x+1)<0$ for $x$ near to 0 with $x \neq 0$.

[^2]Example 1.2.11. In the set $\mathbb{R}_{+}$of positive real numbers, we define

$$
a \nabla b:=\left\{c \in \mathbb{R}_{+}:|a-b| \leq c \leq a+b\right\}
$$

We have that $\mathbb{R}_{+}$with the usual product and $\nabla$ multivalued sum is a multifield, called (real) triangle multifield [58]. We denote this multifield by $\mathcal{T} \mathbb{R}_{+}$.

Note that $a \nabla 0=\{a\}$ and $a \nabla a=\left\{x \in \mathbb{R}_{+}:|x| \leq a\right\}$.
We have some different ways to generalize this construction. If $(F, \leq)$ is an ordered field, we define the triangle multifield $\mathcal{T} F=\left(F_{+}, \nabla, \cdot, 0,1\right)$, by the same prescription,

$$
a \nabla b=\left\{c \in F_{+}:|a-b| \leq c \leq a+b\right\} .
$$

Here, $F_{+}=\{a \in F: a \geq 0\}$. If $(R, P)$ is an ordered ring with supp $(P)=\{0\}$ (for example, $\mathbb{Z}$ ), we define the triangle multiring $\mathcal{T} R=\left(R_{+}, \nabla, \cdot, 0,1\right)$,

$$
a \nabla b=\left\{c \in R_{+}:|a-b| \leq c \leq a+b\right\} .
$$

Again, $R_{+}=\{x \in R: x \geq 0\}$.
kaleid
Example 1.2.12 (Kaleidoscope). Let $n \in \mathbb{N}$ and define $X_{n}=\{-n, \ldots, 0, \ldots, n\} \subseteq \mathbb{Z}$. We define the $n$-kaleidoscope multiring by $\left(X_{n},+, \cdot,-, 0,1\right)$, where $-: X_{n} \rightarrow X_{n}$ is restriction of the opposite map in $\mathbb{Z},+: X_{n} \times X_{n} \rightarrow \mathcal{P}\left(X_{n}\right) \backslash\{\emptyset\}$ is given by the rules:

$$
a+b=\left\{\begin{array}{l}
\{a\}, \text { if } b \neq-a \text { and }|b| \leq|a| \\
\{b\}, \text { if } b \neq-a \text { and }|a| \leq|b|, \\
\{-a, \ldots, 0, \ldots, a\} \text { if } b=-a
\end{array},\right.
$$

and $\cdot: X_{n} \times X_{n} \rightarrow X_{n}$ is is given by the rules:

$$
a \cdot b=\left\{\begin{array}{l}
\operatorname{sgn}(a b) \max \{|a|,|b|\} \text { if } a, b \neq 0 \\
0 \text { if } a=0 \text { or } b=0
\end{array} .\right.
$$

. In this sense, $X_{0}=\{0\}$ and $X_{1}=\{-1,0,1\}=Q_{2}$. For $X_{2}$, we have the following "multioperation" table for the sum:

| + | -2 | -1 | 0 | 1 | 2 |
| :--- | :--- | :--- | :--- | :--- | :--- |
| -2 | $\{-2\}$ | $\{-2\}$ | $\{-2\}$ | $\{-2\}$ | $\{-2,-1,0,1,2\}$ |
| -1 | $\{-2\}$ | $\{-1\}$ | $\{-1\}$ | $\{-1,0,1\}$ | $\{2\}$ |
| 0 | $\{-2\}$ | $\{-1\}$ | $\{0\}$ | $\{1\}$ | $\{2\}$ |
| 1 | $\{-2\}$ | $\{-1,0,1\}$ | $\{1\}$ | $\{1\}$ | $\{2\}$ |
| 2 | $\{-2,-1,0,1,2\}$ | $\{2\}$ | $\{2\}$ | $\{2\}$ | $\{2\}$ |

and the following operation table for the product:

| $\cdot$ | -2 | -1 | 0 | 1 | 2 |
| :--- | :--- | :--- | :--- | :--- | :--- |
| -2 | 2 | 2 | 0 | -2 | -2 |
| -1 | 2 | 1 | 0 | -1 | -2 |
| 0 | 0 | 0 | 0 | 0 | 0 |
| 1 | -2 | -1 | 0 | 1 | 2 |
| 2 | -2 | -2 | 0 | 2 | 2 |

Clearly $\left(X_{n}, \cdot, 1\right)$ is a commutative monoid and $a \cdot 0=0$ for all $a \in X_{n}$.
Now, we will verify that $\left(X_{n},+, \cdot,-, 0,1\right)$ is a multiring.
$i$ - By construction, $a+b=b+a, a+0=\{a\}$ and $0 \in a-a$ for all $a, b \in X_{n}$.
ii $-d \in a+b \Leftrightarrow b \in d-a$ : We divide the proof in cases. Let $a \neq-b$ and suppose without loss of generality that $|a|<|b|$. Thus $a+b=\{b\}$. Hence $d \in a+b$ implies $d=b$. So $b \in b-a=\{b\}$. By symmetry, the same proof applies to the implication $b \in d-a \Rightarrow d \in a+b$. The case $|b|=|a|$ is immediate.
iii $-(a+b)+c=a+(b+c)$ : Again we divide in cases. We suppose without loss of generality that $a, b, c \neq 0$. If $a \neq-b, b \neq-c$, and $|a| \leq|b| \leq|c|$,

$$
(a+b)+c=a+(b+c)=\{c\} .
$$

Similarly, $(a+b)+c=a+(b+c)$ for the cases $|a| \leq|c| \leq|b|,|b| \leq|a| \leq|c|,|b| \leq|c| \leq|a|$, $|c| \leq|a| \leq|b|$ and $|c| \leq|b| \leq|a|$ (under the hypothesis $a \neq-b, b \neq-c$ ).

Now let $a=-b$. We want to prove that $(a-a)+c=a+(-a+c)$. If $|a| \leq|c|$,

$$
(a-a)+c=X_{a}+c=\{c\} \text { and } a+(-a+c)=a+c=\{c\}
$$

If $|c|<|a|$, then

$$
(a-a)+c=X_{a}+c=X_{a} \text { and } a+(-a+c)=a-a=X_{a}
$$

The case $b=-c$ is analogous.
$i v-d(a+b) \subseteq d a+d b$ : If $d=0$ there is nothing to prove. Let $d \neq 0$. If $a \neq-b$, suppose without loss of generality that $|a|<|b|$. Then $a+b=\{b\}$ and $d(a+b)=\{d b\}=d b+d b$.

Now let $a=-b$. We have two cases:
(a) $|d| \leq|a|:$ since $d a=\operatorname{sgn}(d a)|a|$, we have $d a-d a=X_{d a}=X_{a}$ and $d(a-a)=d X_{a} \subseteq X_{a}$.
(b) $|d|>|a|:$ since $d a=\operatorname{sgn}(d a)|d|$, we have $d a-d a=X_{d a}=X_{d}$ and $d(a-a)=d X_{a} \subseteq X_{d}$.

Thus $X_{n}$ is a multiring.

Example 1.2.13 (H-multifield, Example 2.8 in [24]). Let $p \geq 1$ be a prime integer and $H_{p}:=$ $\{0,1, \ldots, p-1\} \subseteq \mathbb{N}$. Now, define the binary multioperation and operation in $H_{p}$ as follow:

$$
\begin{aligned}
& a+b=\left\{\begin{array}{l}
H_{p} \text { if } a=b, a, b \neq 0 \\
\{a, b\} \text { if } a \neq b, a, b \neq 0 \\
\{a\} \text { if } b=0 \\
\{b\} \text { if } a=0
\end{array}\right. \\
& a \cdot b=k \text { where } 0 \leq k<p \text { and } k \equiv a b \bmod p .
\end{aligned}
$$

$\left(H_{p},+, \cdot,-, 0,1\right)$ is a hyperfield such that for all $a \in H_{p},-a=a$. In fact, these $H_{p}$ is a kind of generalization of $K$, in the sense that $H_{2}=K$.

We have to treat sums with some care when we are working with multirings. In order to use the multivalued sum without danger, we define recursively for $n \geq 2$ :

$$
a_{1}+\ldots+a_{n}:=\bigcup_{d \in a_{2}+\ldots+a_{n}} a_{1}+d
$$

In particular, for a multiring $A$, with $a_{1}, \ldots, a_{n} \in A$ and $\sigma \in S_{n}$, we have

$$
a_{1}+a_{2}+\ldots+a_{n}=a_{\sigma(1)}+a_{\sigma(2)}+\ldots+a_{\sigma(n)} .
$$

We also use two conventions: if $Z, W \subseteq R$ and $x \in R, Z+W:=\bigcup\{x+y: x \in Z, y \in W\}$ and $Z+x:=Z+\{x\}=\bigcup\{z+x: z \in Z\}$. We work freely with the immediate consequences of these conventions. For example, from commutativity and associativity is immediate that for all $X, Y, Z \subseteq R, X+Y=Y+X$ and $(X+Y)+Z=Z+(X+Y)$. We return further to these conventions, in the general case of superfields (see for instance Lemma 3.1.5).

Lemma:1.4
Lemma 1.2.14. Let $F$ be a multifield. Then $(a+b) d=a d+b d$ for every $a, b, d \in F$.
Proof. We have $(a+b) d \subseteq a d+b d$ already. For the other inclusion, if $d=0$, it is done. If $d \neq 0$, we have:

$$
\begin{array}{r}
(a d+b d) d^{-1} \subseteq(a d) d^{-1}+(b d) d^{-1}=a d+b d \Rightarrow \\
a d+b d=\left[(a d+b d) d^{-1}\right] d \subseteq(a+b) d .
\end{array}
$$

Then every multifield is in fact a hyperfield, and we use "hyperfield" from now on since it is the prevailing terminology. Now we treat about morphism.
defn:morphism
Definition 1.2.15. Let $A$ and $B$ multirings. A function $f: A \rightarrow B$ is a morphism if for all $a, b, c \in A$ :

$$
\begin{array}{ll}
i-c \in a+b \Rightarrow f(c) \in f(a)+f(b) ; & i v-f(a b)=f(a) f(b) ; \\
i i-f(-a)=-f(a) ; & \\
i i i-f(0)=0 ; & v-f(1)=1 .
\end{array}
$$

The category of multirings with their morphisms will be denoted by MRing.
For multirings, there are types of morphisms that can be considered. Let $f: A \rightarrow B$ a multiring morphism.

- $f$ is a strong morphism if for all $a, b, c \in A$, if $f(c) \in f(a)+f(b)$, then there exist $a^{\prime}, b^{\prime}, c^{\prime} \in A$ with $f\left(a^{\prime}\right)=f(a), f\left(b^{\prime}\right)=f(b), f\left(c^{\prime}\right)=f(c)$ such that $c^{\prime} \in a^{\prime}+b^{\prime}$.
- $f$ is an ideal morphism if for all $a, b, c \in A$, if $f(c) \in f(a)+f(b)$, then exists $c^{\prime} \in A$ with $f\left(c^{\prime}\right)=f(c)$ such that $c^{\prime} \in a+b$. In other words, $f(a+b)=(f(a)+f(b)) \cap \operatorname{Im}(f)$.
- We say that $f$ is a full morphism if it is a strong morphism for all $a, b \in A$ and all $d \in B$,

$$
d \in f(a)+f(b) \Rightarrow \text { exists } c \in a+b \text { such that } d=f(c) .
$$

In other words, $f(a+b)=f(a)+f(b)$.

- We say that $f$ is a strong embedding if $f$ is injective and it is a strong morphism. In this case, $A$ is a submultiring of $B$ if $A \subseteq B$ and the canonical inclusion $\iota: A \hookrightarrow B$ is a strong embedding.
- We say that $f$ is a full embedding if it is a strong embedding and a full morphism ${ }^{3}$.

The different notions of morphisms are related by the following:

$$
\begin{gathered}
\text { Full Morphism } \Rightarrow \text { Ideal Morphism } \Rightarrow \text { Strong Morphism } \\
\text { Full Embedding } \Rightarrow \text { Strong Embedding } \Leftrightarrow \text { Ideal Embedding }
\end{gathered}
$$

The category of hyperfields (respectively multirings) and theirs morphisms will be denoted by $\mathcal{M F}$ (respectively $\mathcal{M R}$ ).

Some of the properties of rings morphisms are not extend to multirings morphisms. Next, are some counterexamples:
ex:2.1

## Example 1.2.16.

$a$ - Let $f: A \rightarrow B$ be a multiring morphism. Define

$$
\operatorname{Ker}(f):=\{a \in A: f(a)=0\} .
$$

$\operatorname{Ker}(f)$ is a submultiring of $A$.
$b-$ Let $f: A \rightarrow B$ be a multiring morphism. If $f$ is injective, them $\operatorname{Im}(f):=\{f(a): a \in A\}$ is embedded in B, but is not a strong embedding and $\operatorname{Im}(f)$ is not a submultiring of $B$ in general. For example, let $R$ be a ring and define a very trivial multioperation * by $a * 0=\{a\}$ for all $a \in R$ and $a * b=R$ if $a, b \neq 0 .(R, *, \cdot, 0,1)$ is a multiring, and considering $R$ as $a$ multiring, the embedding $(R,+, \cdot, 1,0) \hookrightarrow(R, *, \cdot, 0,1)$ is a bijective multiring morphism that is a strong embedding but $(R,+, \cdot, 1,0)$ is not a submultiring of $(R, *, \cdot, 0,1)$. If we consider $K$ as in 1.2.9(b), the inclusion $K \hookrightarrow(R, *, \cdot, 0,1)$ is a multiring morphism that is an embedded and is not a strong embedding.
$c-$ Let $f: \mathbb{R} \rightarrow Q_{2}$ be $f(x)=\operatorname{sgn}(x)$, (with convention that $\operatorname{sgn}(0)=0$ ). $f$ is a multiring morphism, but $f$ is not injective and $\operatorname{Kerf}=\{0\}$. Also $\mathbb{R} / \operatorname{Kerf}$ is not isomorphic to $Q_{2}$.
$d$ - The inclusions functions $Q_{2} \hookrightarrow \mathbb{R}$ and $\mathcal{T} \mathbb{R}_{+} \hookrightarrow \mathbb{R}$ are not multiring morphisms.
$e$ - The inclusion function $\iota: K \rightarrow Q_{2}$ ( $K$ as in 1.2.9(b)) is not a multiring morphism.

### 1.3 Commutative Multialgebra

In the sequel, we extend some terminology of commutative algebra from multirings and hyperfields that could appear throughout this text. Of course, we are not intend to exhaust the theme and for a more detailed exposition, we recommend H. Ribeiro's ph.D Thesis, [23] (in portuguese) or [24].

[^3]Definition 1.3.1 (Definition 2.11 of [24]). An ideal of a multiring $A$ is a non-empty subset of $A$ such that $\mathfrak{a}+\mathfrak{a} \subseteq \mathfrak{a}$ and $A \mathfrak{a}=\mathfrak{a}$. An ideal $\mathfrak{p}$ of $A$ is said to be prime if $1 \notin \mathfrak{p}$ and $a b \in \mathfrak{p} \Rightarrow a \in \mathfrak{p}$ or $b \in \mathfrak{p}$. An ideal $\mathfrak{m}$ is maximal if for all ideals $\mathfrak{a}$ with $\mathfrak{m} \subseteq \mathfrak{a} \subseteq A \Rightarrow$, then $\mathfrak{a}=\mathfrak{m}$ or $\mathfrak{a}=A$. We will denote $\operatorname{Spec}(A)=\{\mathfrak{p} \subseteq A: \mathfrak{p}$ is a prime ideal $\}$.

If $\mathfrak{a}$ is an ideal of $A$, note that $0 \in \mathfrak{a}$ and $-\mathfrak{a} \subseteq \mathfrak{a}$. With the notion of ideal, we define some new multirings structures with the language of commutative algebra in mind:
terminology
Definition 1.3.2 (Definition 2.12 of [24]).
$a-$ If $\left\{A_{i}\right\}_{i \in I}$ is a family of multirings, then the product $\Pi_{i \in I} A_{i}$ is a multiring in the natural (component wise) way.
$b-$ Let $\mathfrak{a} \subseteq A$ be an ideal. Elements of $A / \mathfrak{a}$ are cosets $\bar{a}=a+\mathfrak{a}, a \in A$. More explicitly,

$$
a \equiv b \bmod I \text { if and only if } b \in \bar{a} \text {, if and only if }(b-a) \cap \mathfrak{a} \neq \emptyset .
$$

This is the multialgebra analogous of the usual congruence relation in commutative algebra. We define a multiring structure on $A / \mathfrak{a}$ by $\bar{a}+\bar{b}=\{\bar{c}: c \in a+b\},-\bar{a}=\overline{-a}$, the zero and the unit element of $A / \mathfrak{a}$ are $0=\overline{0}$ and $1=\overline{1}$ respectively and multiplication on $A / \mathfrak{a}$ is defined by $\bar{a} \bar{b}=\overline{a b}$. Note that if $\bar{c} \in \bar{a}+\bar{b}$, then exists $c^{\prime} \in a+b$ such that $\overline{c^{\prime}}=\bar{c}$. The natural arrow $\pi: A \rightarrow A / \mathfrak{a}$ is a strong morphism and as in the ring case it is easily proved that given another multiring morphism $f: A \rightarrow B$ with $f(\mathfrak{a})=\{0\}$, there is a unique morphism $\bar{f}: A / \mathfrak{a} \rightarrow B$ such that $f=\bar{f} \circ \pi$.
$c$ - Let $S$ be a multiplicative set in $A$. Elements of $S^{-1} A$ have the form $a / s, a \in A, s \in S$, $a / s=b / t$ if and only if atu $=b s u$ for some $u \in S .0=0 / 1,1=1 / 1$ and the operations are defined by $(a / s) \cdot(b / t)=a b / s t$, and $c / v \in a / s+b / t$ if and only if cstv $\in a t u v+b s u v$ for some $v \in S$. The natural arrow $\rho: A \hookrightarrow S^{-1} A$ is a strong morphism and given a multirng morphism $f: A \rightarrow B$ with $f(S) \subseteq B^{\times}$, then exists a unique morphism $\bar{f}: S^{-1} A \rightarrow B$ such that $f=\bar{f} \circ \rho$.
$d$ - If $D$ is a multidomain, we define the multifield of fractions ff $(D):=(D \backslash\{0\})^{-1} D$.
Let $X$ be a subset of a multiring $A$. We define the ideal generated by $X$ by

$$
\langle X\rangle:=\bigcap\{\mathfrak{a} \subseteq A: X \subseteq \mathfrak{a}, \mathfrak{a} \text { is an ideal }\}
$$

If $X \neq \emptyset$, we have that $\langle X\rangle=\bigcup\left\{\lambda_{1} x_{1}+\cdots+\lambda_{n} x_{n}: n \geq 1, \lambda_{i} \in A, x_{i} \in X\right.$, for all $\left.i=1, \ldots, n\right\}$. In particular

$$
\langle a\rangle=\sum A a:=\left\{\sum_{j=1}^{n} \lambda_{j} a: \lambda_{1}, \ldots, \lambda_{n} \in A, n \geq 1 .\right\} .
$$

If $A$ is a hyperring then $\sum A a=A a$.
lem:iso
Proposition 1.3.3 (Proposition 2.13 of [24). Let $A$ and $B$ be multirings and $\varphi: A \rightarrow B$ a surjective morphism. Consider $\bar{\varphi}: A / \operatorname{Ker}(\varphi) \rightarrow B$ the induced morphism. Then the following are equivalent:
$i-\varphi$ is a strong morphism and if $\varphi(a)=\varphi\left(a^{\prime}\right)$ for $a, a^{\prime} \in A$, then $\left(a-a^{\prime}\right) \cap \operatorname{ker}(\varphi) \neq \emptyset$.
ii- $\varphi$ is an ideal morphism.
iiii- $\bar{\varphi}$ is an isomorphism.
Proof. $i) \Rightarrow i i)$ : Assume that $\varphi(a) \in \varphi(b)+\varphi(c)$. Since $\varphi$ is a strong morphism, exists $a^{\prime}, b^{\prime}, c^{\prime} \in A$ with $\varphi\left(a^{\prime}\right)=\varphi(a), \varphi\left(b^{\prime}\right)=\varphi(b), \varphi\left(c^{\prime}\right)=\varphi(c)$ such that $a^{\prime} \in b^{\prime}+c^{\prime}$. By hypothesis, exists $b^{\prime} \in b+i$ and $c^{\prime} \in c+j$ such that $i, j \in \operatorname{ker}(\varphi)$. Then $a^{\prime} \in b^{\prime}+c^{\prime} \subseteq(b+c)+(i+j)$ and so exists $x \in i+j \subseteq \operatorname{ker}(\varphi)$ such that $a^{\prime} \in b+c+x$. Thus exist $a^{\prime \prime} \in a^{\prime}-x$ with $a^{\prime \prime} \in b+c$ and note that $\varphi\left(a^{\prime \prime}\right)=\varphi\left(a^{\prime}\right)=\varphi(a)$.
$i i) \Rightarrow$ iii $)$ : Let $a, b \in A$ such that $\varphi(a)=\bar{\varphi}(\bar{a})=\bar{\varphi}(\bar{b})=\varphi(b)$. By hypothesis exist $x \in a-b$ such that $x \in \operatorname{ker}(\varphi)$ and so $\bar{a}=\bar{b}$ in $A / \operatorname{ker}(\varphi)$, proving the injectivity of $\bar{\varphi}$. Since $\varphi$ is a strong morphism, if $\bar{\varphi}(\bar{a}) \in \bar{\varphi}(\bar{b})+\bar{\varphi}(\bar{c})$, then exists $a^{\prime}, b^{\prime}, c^{\prime} \in A$ with $\varphi\left(a^{\prime}\right)=\varphi(a), \varphi\left(b^{\prime}\right)=\varphi(b), \varphi\left(c^{\prime}\right)=\varphi(c)$ such that $a^{\prime} \in b^{\prime}+c^{\prime}$. By hypothesis, it is easy to see that $\overline{a^{\prime}}=\bar{a}, \overline{b^{\prime}}=\bar{b}, \overline{c^{\prime}}=\bar{c}$ and so $\bar{a} \in \bar{b}+\bar{c}$ in $A / \operatorname{ker}(\varphi)$. Thus $\bar{\varphi}$ is an isomorphism.
iii $) \Rightarrow i)$ : Assume that $\varphi(a)=\varphi\left(a^{\prime}\right)$ for $a, a^{\prime} \in A$. Then $\bar{\varphi}(\bar{a})=\bar{\varphi}\left(\bar{a}^{\prime}\right)$ and hence $\bar{a}=\overline{a^{\prime}}$, which means that $\left(a-a^{\prime}\right) \cap \operatorname{ker}(\varphi) \neq \emptyset$. Therefore $0=\varphi(0) \in \varphi(a)-\varphi\left(a^{\prime}\right)$, and by hypothesis there exist $i \in a-a^{\prime}$ such that $\varphi(i)=\varphi(0)=0$. On the other hand, we have $\varphi=\bar{\varphi} \circ \pi$, where $\pi: A \rightarrow A / \operatorname{ker}(\varphi)$. Then $\varphi$ is a composition of strong morphisms and so $\varphi$ is strong itself.
teo:iso
Theorem 1.3.4 (Isomorphism Theorem, 2.14 of [24]). Let $A$ and $B$ be multirings and $\varphi: A \rightarrow B$ an ideal morphism. Then $\operatorname{Im}(\varphi)$ is a multiring (contained in $B$ ) with the structure induced by the domain $A$, and the induced morphism $\bar{\varphi}: A / \operatorname{Ker}(\varphi) \rightarrow \operatorname{Im}(\varphi)$ is an isomorphism.
Proof. By the previous Proposition, it is enough to prove that $\operatorname{Im}(\varphi)$ is a multiring and this is accomplished by proving the associativity property for $\operatorname{Im}(\varphi)$. Assume that $\varphi(x) \in \varphi(p)+\varphi(w)$ with $\varphi(p) \in \varphi(u)+\varphi(v)$. Since $\varphi$ is an ideal morphism, exists $x^{\prime} \in p+w$ and $p^{\prime} \in u+v$ such that $\varphi\left(x^{\prime}\right)=\varphi(x)$ and $\varphi\left(p^{\prime}\right)=\varphi(p)$. Then, by the same argument as the previous Lemma, it should exist $i \in \operatorname{Ker}(\varphi)$ such that $p \in i+p^{\prime}$. Then $p \in i+(u+v) \subseteq(i+u)+v$ and thus exist $u^{\prime} \in i+u$ such that $p \in u^{\prime}+v$. Then exist $q \in v+w$ with $x \in u^{\prime}+q$. Thus $\varphi(x) \in \varphi\left(u^{\prime}\right)+\varphi(q)=\varphi(u)+\varphi(q)$ and $\varphi(q) \in \varphi(v)+\varphi(w)$.

Lemma:1.1
Lemma 1.3.5 (Lemma 2.16 of [24). Let $A$ be a multiring. Then:
$a-$ an ideal $\mathfrak{p}$ of $A$ is prime if and only if $A / \mathfrak{p}$ is a multidomain.
$b-$ An ideal $\mathfrak{m}$ is maximal if and only if for all $a \neq 0$ in $A / \mathfrak{m}$, exists $t_{1}, \ldots, t_{n}$ such that

$$
1 \in a t_{1}+\cdots+a t_{n} .
$$

In particular, maximal ideals are prime and if $A$ is a hyperring, an ideal $\mathfrak{m}$ is maximal if and only if $A / \mathfrak{m}$ is a multifield.
Proof.
a - The same of the ring case.
$\mathrm{b}-\Rightarrow$ : Let $\bar{a} \in A / \mathfrak{m}$ non-zero, that is, $a \notin \mathfrak{m}$. Since $m$ is maximal, the ideal generated by $\mathfrak{m} \cup\{a\}$, namely

$$
I=\bigcup\left\{m+a t_{1}+\cdots+a t_{n}: n \geq 1 \text { and } t_{i} \in A\right\}
$$

is improper. Then exists $m \in \mathfrak{m}$ and $t_{1}, \cdots, t_{n} \in A$ such that $1 \in m+a t_{1}+\cdots+a t_{n}$ and so $\overline{1} \in \bar{a} \overline{t_{1}}+\cdots+\bar{a} \overline{t_{n}}$.
$\Leftarrow$ Let $a \notin \mathfrak{m}$. By the property valid in $A / \mathfrak{m}$, exists $m \in \mathfrak{m}$ and $t_{1}, \cdots, t_{n} \in A$ such that $1 \in m+a t_{1}+\cdots+a t_{n}$ and so the ideal generated by $\mathfrak{m} \cup\{a\}$ is improper. Then $\mathfrak{m}$ is maximal.

Proposition 1.3.6 (Proposition 2.17 of [24]).
$a$ - Let $A$ be a multiring, $I \subseteq A$ an ideal and $S \subseteq A$ be a multiplicative subset of $A$. Then $(S / I)^{-1} A / I \cong S^{-1} A / S^{-1} I$.
$b$ - Let $\left\{A_{i}\right\}_{i \in I}$ be a family of multirings and $\mathfrak{a}_{i} \subseteq A_{i}$ be an ideal of $A_{i}$ for every $i \in I$. Then

$$
\prod_{i \in I} A_{i} / \mathfrak{a}_{i} \cong \prod_{i \in I} A_{i} / \prod_{i \in I} \mathfrak{a}_{i} .
$$

Proof. For the item (a), consider the morphism $f: S^{-1} A \rightarrow(S / I)^{-1} A / I$ given by $f(a / s)=\bar{a} / \bar{s}$ and apply the Theorem 1.3.4. For the item (b), the same strategy holds with the morphism $g: \prod_{i \in I} A_{i} \rightarrow \prod_{i \in I} A_{i} / \prod_{i \in I} \mathfrak{a}_{i}$ given by $g\left(a_{i}\right)_{i \in I}=\left(\bar{a}_{i}\right)_{i \in I}$.

Now, we present the main construction related to quadratic forms, that we baptize "Marshall's quotient". This kind of quotient appears naturally in the context of abstract theories of quadratic forms, as we will have the opportunity to see later in the text.
defn:strangeloc
Definition 1.3.7 (Example 2.6 of [47]). Fix a multiring $A$ and a multiplicative subset $S$ of $A$. Define an equivalence relation $\sim$ on $A$ by $a \sim b$ iff as $=b t$ for some $s, t \in S$. Denote by $\bar{a}$ the equivalence class of $a$ and set $A / m S=\{\bar{a}: a \in A\}$. Defining $-\bar{a}=\overline{-a}, \bar{a} \bar{b}=\overline{a b}$ and

$$
\bar{a}+\bar{b}=\{\bar{c}: c v \in a s+b t, \text { for some } s, t, v \in S\},
$$

we have that $\left(A /{ }_{m} S,+, \cdot,-, \overline{0}, \overline{1}\right)$ is a multiring, called the Marshall's quotient of $A$ by $S$.
Let $S$ be a non-empty subset of a multiring $A$. We define the ideal generated by $S$ by $\langle S\rangle:=\bigcap\{\mathfrak{a} \subseteq A$ ideal $: S \subseteq \mathfrak{a}\}$. If $S=\left\{a_{1}, \ldots, a_{n}\right\}$, we have that

$$
\left\langle a_{1}, \ldots, a_{n}\right\rangle=\sum A a_{1}+\ldots+\sum A_{n}, \text { where } \sum A a=\bigcup_{n \geq 1}\{\underbrace{a+\ldots+a}_{n \text { times }}\} .
$$

If $A$ is a hyperring, then $\sum A a=A a$.
Proposition 1.3.8 (Proposition 2.19 of [24]). Let $A, B$ be a multiring and $S \subseteq A$ a multiplicative subset of $A$. Then for every morphism $f: A \rightarrow B$ such that $f[S]=\{1\}$, there exist a unique morphism $\tilde{f}: A /{ }_{m} S \rightarrow B$ such that the following diagram commute:

where $\pi: A \rightarrow A /{ }_{m} S$ is the canonical projection $\pi(a)=\bar{a}$.
Proposition 1.3.9 (Proposition 2.20 of [24]). Let $A$ be a multiring, $I \subseteq A$ an ideal and $S \subseteq A$ a multiplicative subset such that $I \subseteq S$. Define $S / I=\{\bar{s}: s \in S\}$ (modulo $I$ ). Then

$$
(A / I) / m(S / I) \cong A /{ }_{m} S .
$$

Proof. Define $\Phi: A / I \rightarrow A /{ }_{m} S$ given by $\Phi\left(\bar{a}^{I}\right)=\bar{a}^{S}$ and use the previous Proposition.

Proposition 1.3.10 (Proposition 2.21 of [24]). Let $A$ be a multiring and $P, S \subseteq A$ multiplicative subsets of $A$ such that $P \subseteq S$. Then

$$
A /{ }_{m} S \cong P^{-1} A /{ }_{m} P^{-1} S
$$

### 1.4 Ordering Structures and Artin-Schreier Theorem

 defn:mforderingDefinition 1.4.1 (Page 8 of 47]). Let $F$ be a hyperfield. A subset $P$ of $F$ is called an ordering if $P+P=\subseteq P, P \cdot P \subseteq P, P \cup-P=F$ and $P \cap-P=\{0\}$. The real spectrum of a hyperfield $F$, denoted $\operatorname{Sper}(F)$, is defined to be the set of all orderings of $F$.

Definition 1.4.2 (Page 8 of [47]). A preordering of a hyperfield $F$ is defined to be a subset $T$ of $F$ satisfying $T+T \subseteq T, T \cdot T \subseteq T$ and $F^{2} \subseteq T$. Here, $F^{2}:=\left\{a^{2}: a \in F\right\}$. A hyperfield $F$ is said to be real if $-1 \notin \sum F^{2}$. If $F$ is real, then $-1 \neq 1$. A preordering $T$ of $F$ is said to be proper if $-1 \notin T$.
lem:3.2marshall
Lemma 1.4.3 (Lemma 3.2 of [47]). Suppose $F$ is a hyperfield with $-1 \neq 1$. For a preordering $T$ of $F$, the following are equivalent:

$$
i-T \text { is proper. }
$$

$i i-T \neq F$.
Proof. $(i) \Rightarrow(i i)$ is just the definition. For $(i i) \Rightarrow(i)$, suppose that $-1 \in T$ and let $a \in F$. If $a=0$ then $a \in T$. Suppose $a \neq 0$. Fix $b \in 1+a$. Then $b^{2} \in 1+a+a+a^{2}$, so $b^{2} \in 1+u+a^{2}, u \in a+a$. Then $u \in b^{2}-1-a^{2} \in T . u / a \in 1+1$, so $u / a \in T$. Since $-1 \neq 1, u \neq 0$ and $\dot{T}$ is a subgroup of $\dot{F}$, then $a / u=(u / a)^{-1} \in T$. Hence $a=(a / u) u \in T$.
lem:3.3marshall
Lemma 1.4.4 (Lemma 3.3 of [47).
$a-A$ preordering which is maximal and proper is an ordering.
$b-F$ has ordering if and only if $F$ is real.
Proof. a - Let $P$ be a preordering of the hyperfield $F$ which is maximal and proper. If $a \in F$, then $P-a P$ is also a preordering. If $-1 \in P-a P$, then there exists $s, t \in P$ such that $-1 \in s-a t$. If $t=0$, then $-1=s \in P$, a contradiction. Thus $t \neq 0$. Then $a t \in 1+s$, so $a \in 1 / t+s / t \subseteq P$. If $-1 \notin P-a P$, then by maximality of $P,-a \in P$. This proves that $P \cup-P=F$. If $s \in P \cap-P$, $s \neq 0$, then $s=-t \in P$, so $-1=s / t \in P$, contradiction. This proves that $P \cap-P=\{0\}$.
b-By Zorn's Lemma, every preordering is containing in an ordering. This fact with the item (a) proves the desired.

For a preordering $T$ of $F$, we denote by $X_{T}$ the set of all orderings of $F$ with $T \subseteq F$.
prop:3.4marshall
Proposition 1.4.5 (Generalized Artin-Schreier Theorem (Proposition 3.4 of [47)). Let $F$ be an hyperfield and $T$ a proper preordering of $F$. Then $T=\bigcap_{P \in X_{T}} P$, where $X_{T}=\{P \in \operatorname{Sper}(F): T \subseteq$ $P\}$.

Proof. The inclusion " $\subseteq$ " is immediate. For the inclusion " $\supseteq$ ", fix $a \in F, a \notin T$. Then $T-a T$ is a proper preordering of $F$ (the argument is the same of 1.4.4). By the Zorn's Lemma, there exists a maximal and proper preordering $P$ such that $T-a T \subseteq P$. By 1.4.4, $P$ is an ordering, and $-a \in P$, so $a \notin P$.

### 1.5 Real Reduced hyperfields

Consider the hyperfield $Q_{2}$. $\{0,1\}$ is an ordering on $Q_{2}$. For any ordering $P$ on a hyperfield $F, Q_{P}(F)=F / m P \cong Q_{2}$ by a unique isomorphism. Orderings of a hyperfield $F$ correspond bijectively to a multiring homomorphism $\sigma: F \rightarrow Q_{2}$ via $P=\sigma^{-1}(\{0,1\})$.
prop:4.1marshall
Proposition 1.5.1 (Proposition 4.1 of [47]). For a real hyperfield $F$ are equivalent:
$a$ - The multiring morphism $F \rightarrow Q_{\text {red }}(F)$ is an isomorphism;
$b-\sum F^{2}=\{0,1\} ;$
$c$ - For all $a \in F, a^{3}=a$ and $(a \in 1+1) \Rightarrow(a=1)$.
Proof. (a) $\Leftrightarrow(\mathrm{b})$ Is just the general fact that if $\sigma: F \rightarrow K$ is a morphism of real hyperfields, then $\sigma\left(\sum F^{2}\right) \subseteq \sum K^{2}$ and that $\sum Q_{r e d}(F)^{2}=\{0,1\}$.
$(\mathrm{a}) \Rightarrow(\mathrm{c}) Q_{\text {red }}(F)$ already satisfy $a^{3}=a$ for all $a$ and $1+1=\{1\}$.
(c) $\Rightarrow$ (b) We have $a^{2}=1$ for all $a \neq 0$ and $\underbrace{1+1+\ldots+1}_{n}=\{1\}$ by induction on $n$. It follows that $\sum F^{2}=F^{2}=\{0,1\}$.
defn:mfrealreduced
Definition 1.5.2. A hyperfield $F$ is said to be real reduced if satisfies the equivalent conditions of Proposition 1.5.1.

A morphism of real reduced hyperfield is just a morphism of hyperfields. The category of real reduced hyperfields will be denoted by $\mathcal{M} \mathcal{F}_{\text {red }}$.
cor:4.2marshall
Corollary 1.5.3 (Corollary 4.2 of [47]). A hyperfield $F$ is real reduced if and only if $a^{3}=a$ for all $a \in F$ and $a \in 1+1 \Rightarrow a=1$.

Proof. $(\Rightarrow)$ is already done. For $(\Leftarrow)$, by Proposition 1.5 .1 is suffice to prove that $F$ is real. Therefore, suppose that $a^{3}=a$ for all $a \in F$ and $a \in 1+1 \Rightarrow a=1$. Then $\sum F^{2}=\{0,1\}$. If $-1 \in\{0,1\}$, then $-1=0$, so $1=0$ or $-1=1$, so $0 \in 1+1=\{1\}$. In both cases, we conclude that $1=0$, contradiction. Thus $-1 \notin \sum F^{2}$, then $F$ is real.

For any proper preordering $T$ of a real reduced hyperfield $F, Q_{T}(F)$ is a real reduced hyperfield. In particular, $Q_{r e d}(F)$ is a real reduced hyperfield. If $p: F_{1} \rightarrow F_{2}$ is a multiring homomorphism of real hyperfields, then $p$ induces a morphism $Q_{r e d}\left(F_{1}\right) \rightarrow Q_{r e d}\left(F_{2}\right)$. In this way, $Q_{r e d}$ defines a functor (a reflection) from the category of real hyperfields onto the subcategory of real reduced hyperfields.
lem:4.3marshall
Proposition 1.5.4 (Lemma 4.3 of [47]). Let $F$ be a real reduced hyperfield, $T=\sum F^{2}$. For any $a, b \in \dot{F}$,

$$
(a+b)^{*}=(T a+T b)^{*}=\{c \in \dot{F}: \forall \sigma \in \operatorname{Sper}(F), \sigma(c)=\sigma(a), \text { or } \sigma(c)=\sigma(b)\} .
$$

Proof. Since $F$ is a real reduced hyperfield, $T=\{0,1\}$, so $T a+T b=\{0, a, b\} \cup(a+b)$. In particular, $F=T-T=\{0,1,-1\} \cup(1-1)$. To prove $(a+b)^{*}=(T a+T b)^{*}$, it remains to show $a, b \in a+b$. By symmetry, it suffices to show $a \in a+b$. If $a \neq \pm b$, then $b / a \neq \pm 1$ so $b / a \in 1-1$, i.e, $b \in a-a$ and so $a \in a+b$. If $a=b, 1 \in 1+1 \Rightarrow a \in a+a=a+b$, and if $a=-b,-b \in-b-b \Rightarrow a \in a-b \Rightarrow a \in a+b$. Therefore $(a+b)^{*}=(T a+T b)^{*}$.

If $c \in T a+T b$, then $\sigma(a)=\sigma(b)$ implies that $\sigma(c)=\sigma(a)$. Thus $\sigma(c)=\sigma(a)$ or $\sigma(c)=\sigma(b)$ for any $\sigma \in \operatorname{Sper}(F)$. Conversely suppose this holds for any $\sigma$. Then $\sigma(b / a)=1$ implies $\sigma(c / a)=1$ for any $\sigma$, so by Proposition 1.4.5, $c / a \in T+T(b / a)$. Multiplying by $a$, this yields $c \in T a+T b$ as required.

Real reduced hyperfields have a natural representation in terms of functions:
cor:4.4marshall
Theorem 1.5.5 (Local-Global principle, Corollary 4.4 of [47]). For any real reduced hyperfield $F$, the natural embedding $F \hookrightarrow Q_{2}^{\operatorname{Sper}(F)}$ is a strong embedding.
Proof. Let $F$ be a real reduced hyperfield and $T=\sum F^{2}=\{0,1\}$. By Proposition 1.4.5,

$$
\{0,1\}=\bigcap_{P \in X_{T}} P
$$

or in other words, 1 is the unique element that is positive in all orderings. Hence, if $\sigma(a)=\sigma(b)$ for all $\sigma \in X_{T}$, then $a b$ is positive in all orderings, so $a b=1$ and as $a^{2}=1$, we have $a=b$. Therefore, the multiring morphism from $F$ to $Q_{2}^{\operatorname{Sper}(F)}$ defined by $a \mapsto(\sigma(a))_{\sigma \in \operatorname{Sper}(F)}$ is injective.

It remains to show that if $\sigma(c) \in \sigma(a)+\sigma(b)$ for all $\sigma \in \operatorname{Sper}(F)$ then $c \in a+b$. If $a=0$, then $\sigma(c)=\sigma(b)$ for all $\sigma \in X_{T}$, so by the argument above. $b=c$. Similarly, if $b=0$ then $c=a$ and if $c=0$, then $b=-a$. Suppose now that $a, b, c$ are not zero. Then $c \in a+b$ by Proposition 1.5.4.

In particular, for any real reduced hyperfield, $\operatorname{Sper}(F)$ separate points of $F$ and $c \in a+b \subseteq F$ if and only if, for every $\sigma: F \rightarrow Q_{2}, \sigma(c) \in \sigma(a)+\sigma(b)$.

### 1.6 The Positivstellensatz

Let $A$ be a multiring. A subset $P$ of $A$ is an ordering if $P+P \subseteq P, P P \subseteq P, P \cup-P=A$ and $P \cap-P$ is a prime ideal of $A$ (called the support of $A$ ). Orderings of a multiring $A$ correspond bijectively to multiring homomorphisms $\sigma: A \rightarrow Q_{2}$ via $P=\sigma^{-1}(\{0,1\})$. For a prime ideal $\mathfrak{p}$ of $A$, orderings on $A$ having support contained in $\mathfrak{p}$ (resp., containing $\mathfrak{p}$, resp., equal to $\mathfrak{p}$ ) correspond bijectively to orderings on the localization of $A$ (resp., on $A / \mathfrak{p}$, on $f f(A / \mathfrak{p})$ ). The real spectrum of $A$, denoted $\operatorname{Sper}(A)$, is the set of all orderings of $A$.

A preordering of a multiring $A$ is a subset $T$ of $A$ satisfying $T+T \subseteq T, T T \subseteq T$ and $A^{2} \subseteq T$. A preordering $T$ of $A$ is said to be proper if $-1 \notin T$. Every ordering is a proper preordering. $\sum A^{2}$ us a preordering, and is the unique smallest preordering of $A$. A multiring $A$ is said to be semi real if $-1 \notin \sum A^{2}$.

Fix a preordering $T$ of $A$. Define $X_{T}:=\{\sigma \in \operatorname{Sper}(A): \sigma(T)=\{0,1\}\}$. A $T$-module in $A$ is defined to be a subset $M$ of $A$ satisfying $M+M \subseteq M, T M \subseteq M$, and $1 \in M$ (so $T \subseteq M$ ).
prop:5.2marshall
Proposition 1.6.1 (Proposition 5.2 of [47]). Suppose $T$ is a preordering of $A$ and $M$ is a $T$ module in $A$ which is maximal subject to $-1 \notin M$. Then $M \cap(-M)$ is a prime ideal of $A$, and $M \cup(-M)=A$.

Proof. First we show that $\mathfrak{p}=M \cap-M$ is an ideal. Let $M^{\prime}=\{a \in A:(a+a) \cap M \neq \emptyset\}$. Then $M^{\prime} \supseteq M$ and $M^{\prime}$ is a $T$-module. If $-1 \in M^{\prime}$, then $(-1-1) \cap M \neq \emptyset$, say $a \in(-1-1) \cap M$. Then $-1 \in 1+a \subseteq M$, a contradiction. Thus $-1 \notin M^{\prime}$. By maximality of $M, M=M^{\prime}$. By construction, we have $\mathfrak{p}+\mathfrak{p} \subseteq \mathfrak{p},-\mathfrak{p}=\mathfrak{p}$ and $T \mathfrak{p} \subseteq \mathfrak{p}$. Suppose $a \in A, b \in \mathfrak{p}$ are given. Fix $c \in 1+a$. Then $c^{2} \in 1+a+a+a^{2}$, so $c^{2} \in 1+d+a^{2}$ for some $d \in a+a$. Then $d \in c^{2}-1-a^{2}$, so $d b \in c^{2} b-b-a^{2} b \subseteq \mathfrak{p} \subseteq M$. At same time, $d b \in(a+a) b \subseteq a b+a b$. This proves $a b \in M^{\prime}=M$. A similar argument shows that $a b \in-M$. Thus $a b \in M \cap-M=\mathfrak{p}$. This proves that $\mathfrak{p}$ is an ideal of $A$.

Next we show that $\mathfrak{p}$ is prime. Suppose $a b \in \mathfrak{p}, a \notin \mathfrak{p}, b \notin \mathfrak{p}$. Replacing $a$ by $-a$ and $b$ by $-b$ if necessary, we can assume $a \notin M, b \notin M$. Thus -1 lies in the $T$-module $M+\sum a T$ and also in the $T$-module $M+\sum b T$. Then $-b^{2} \in M b^{2}+\sum a b^{2} T \subseteq M$ (using the fact that $a b \in \mathfrak{p}$ ), so $b^{2} \in \mathfrak{p}$. Writing $-1 \in q+c, q \in M, c \in \sum b t_{i}, t_{i} \in T$, we have $-c \in 1+q$, so $c^{2} \in 1+q+q+q^{2}$. on the other hand, $c^{2} \in \sum b^{2} t_{i} t_{j} \subseteq \mathfrak{p}$. This implies $-1 \in-c^{2}+q+q+q^{2} \subseteq M$, a contradiction. This proves that $\mathfrak{p}$ is a prime ideal.

Finally, we prove that $A=M \cup-M$. Suppose $a \in A$ with $a \notin M$ and $a \notin-M$. Then $-1 \in M+\sum a T$ and $-1 \in M-\sum a T$. Multiplying by $a^{2}$, and noting that $a\left(\sum a T\right) \subseteq T$, this yields $-a^{2} \in M+t_{1} a-a^{2}$ and $-a^{2} \in M-t_{2} a$, for some $t_{1}, t_{2} \in T$. Then $-t_{1} a \in a^{2}+M \subseteq M$, and $t_{2} a \in a^{2}+M \subseteq M$, so $t_{1} t_{2} a \in \mathfrak{p}$. This is not possible. If either of $t_{1}$ or $t_{2}$ is in $\mathfrak{p}$, then $-a^{2} \in M$, so $-1 \in M+\sum a T \Rightarrow a \in-M+\sum\left(-a^{2}\right) T$, and $-1 \in M-\sum a T \Rightarrow-a \in M+\sum\left(-a^{2}\right) T$, then $a \in \mathfrak{p}$. If $a \in \mathfrak{p}$, then $a \in M$ (and also $a \in-M$ ), which contradiction our assumption. This proves $A=M \cup-M$.
cor:5.3marshall
Corollary 1.6.2 (Corollary 5.3 of [47]). $\operatorname{Sper}(A) \neq \emptyset$ if and only if $-1 \notin \sum A^{2}$. For a preordering $T$ of $A, X_{T} \neq \emptyset$ if and only if $T$ is proper.

Proof. The first assertion follows from the second. If $X_{T} \neq \emptyset$ then clearly $T$ is proper. Suppose now that $T$ is proper. Use Zorn's Lemma to choose a maximal proper preordering $P$ in $A$ with $T \subseteq P$, and a $P$-module $M$ of $A$ maximal subject to $-1 \notin M$. If $P \neq M$ then for any $a \in M \backslash P$, $P+\sum a P$ is a preordering and $P+\sum a P \subseteq M$, so $P+\sum a P$ is proper. This contradicts the maximality of $P$. It follows that $P=M$. Proposition 1.6.1 implies that $P$ is an ordering.

For a fixed preordering $T$ of $A$ we have a multiring homomorphism $A \rightarrow Q_{2}^{X_{T}}$ (the product multiring), given by $a \mapsto \bar{a}$, where $\bar{a}$ is defined by $\bar{a}(\sigma)=\sigma(a)$ for all $\sigma \in X_{T}$.
prop:5.4marshall
Proposition 1.6.3 (Proposition 5.4 of [47]). Suppose $c, d \in A$. Then $\bar{c} \geq 0 \Rightarrow \bar{d}=0$ holds on $X_{T}$ (i.e, $\sigma(c) \geq 0 \Rightarrow \sigma(d)=0$ ) if and only if $-d^{2 k} \in T+\sum A^{2} c$ for some integer $k \geq 0$.

Proof. $(\Rightarrow)$ Let $B=S^{-1} A, T^{\prime}=S^{-1} T$, where $S:=\left\{d^{2 k}: k \geq 0\right\}$, and consider the $T$-module $T+\sum A^{2} c$ and the $T^{\prime}$-module $T^{\prime}+\sum B^{2} c$. If $-S \cap\left(T+\sum A^{2} c\right)=\emptyset$, then $-1 \notin T^{\prime}+\sum B^{2} c$, so there is a $T^{\prime}$-module $M$ in $B$ containing $T^{\prime}+\sum B^{2} c$ and maximal subject to $-1 \notin M$. By Proposition 1.6.1, $\mathfrak{p}:=M \cap-M$ is a prime ideal. Also, $T^{\prime} \subseteq M$, so $\left(T^{\prime}+\mathfrak{p}\right) \cap\left(-T^{\prime}+\mathfrak{p}\right)=\mathfrak{p}$. It follows that the preordering $T^{\prime \prime}:=\left\{(a+\mathfrak{p}) /(b+\mathfrak{p}): a, b \in T^{\prime}, b \notin p\right\}$ is a proper preordering in the hyperfield $F:=f f(A / \mathfrak{p})$. Since $d \notin \mathfrak{p}(d$ is invertible in $B)$, it follows from our assumption that $c+\mathfrak{p} \in \notin P$ for all orderings $P$ of $F$ containing $T^{\prime \prime}$. According to Proposition 1.4.5, this implies that $c+\mathfrak{p} \in-T^{\prime \prime}$. This yields elements $s, t \in T^{\prime}+\mathfrak{p}$ with $s, t \notin \mathfrak{p}$ such that $-s c=t$. Then $s t \in T^{\prime}+\mathfrak{p} \subseteq M$ and $-s t=s^{2} c \in \sum B^{2} c \subseteq M$, so $s t \in M \cap-M=\mathfrak{p}$, a contradiction.
$(\Leftarrow)$ We already know that $\sigma\left(d^{2 k}\right) \geq 0$ for all $\sigma \in X_{T}$. If $-d^{2 k} \in T+\sum A^{2} c$, then $-\sigma\left(d^{2 k}\right) \geq 0$ for all $\sigma \in X_{T}$. Hence $\sigma\left(d^{2 k}\right)=-\sigma\left(d^{2 k}\right)=0$ for all $\sigma \in X_{T}$, and this implies that $\sigma(d)=0$ for all $\sigma \in X_{T}$.

Corollary 1.6.4 (Corollary 5.5 of [47]).
$a-\bar{a}=0$ on $X_{T}$ if and only if $-a^{2 k} \in T$ for some $k \geq 0$.
$b-\bar{a}=1$ on $X_{T}$ if and only if $-1 \in T-\sum A^{2} a$.
$c-\bar{a} \geq 0$ on $X_{T}$ if and only if $-a^{2 k} \in T-\sum A^{2} a$ for some $k \geq 0$.
$d$ - Fix $a \in b^{2}+c^{2}$. Then $\bar{b}=\bar{c}$ on $X_{T}$ if and only if $-a^{2 k} \in T-\sum A^{2} b c$ for some $k \geq 0$.
Proof. Apply Proposition 1.6 .3 as follows: (a) take $c=0, d=a$. (b) Take $c=-a, d=1$. (c) Take $c=-a, d=a$. (d) Take $c=-b c, d=a$.

### 1.7 Real Ideals

We indicate briefly how the theory of real ideals and real prime ideals extends to multirings. An ideal $\mathfrak{a}$ in a multiring $A$ is said to be real if $\left(\sum a_{i}^{2}\right) \cap \mathfrak{a} \neq \emptyset \Rightarrow a_{i} \in \mathfrak{a}$ for each $i$. Every real ideal is radical in the sense that $a^{2} \in \mathfrak{a} \Rightarrow a \in \mathfrak{a}$, i.e, $\mathfrak{a}$ is the intersection of prime ideals of $A$. The converse is not true.
prop:6.1marshall
Proposition 1.7.1 (Proposition 6.1 of [47]). For a prime ideal $\mathfrak{p}$ in a multiring A, the following are equivalent:
$a-\mathfrak{p}$ is real.
$b$ - The residue hyperfield $f f(A / \mathfrak{p})$ is real.
$c-\mathfrak{p}$ is the support of some ordering of $A$.
Proof. (a) $\Rightarrow$ (b) If $-1+\mathfrak{p} \in \sum a_{i}^{2}+\mathfrak{p}$, then $0 \in 1+\sum a_{i}^{2}+\mathfrak{p}$, and $\left(1+\sum a_{i}^{2}\right) \cap \mathfrak{p} \neq \emptyset$. As $\mathfrak{p}$ is real, $1 \in \mathfrak{p}$, contradiction. Then $-1 \notin \sum(A / \mathfrak{p})^{2}$, and therefore $-1 \notin \sum f f(A / \mathfrak{p})^{2}$.
(b) $\Rightarrow$ (c) By Proposition 1.4.4 $f f(A / \mathfrak{p})$ has an ordering $P$. Let $\tilde{P}=\left\{a_{i}, b_{i}: a_{i} / b_{i} \in P\right\}$ and $Q=q^{-1}[\tilde{P}]$, where $q: A \rightarrow A / \mathfrak{p}$ is the canonical projection. Then $Q$ is the desired ordering.
$(\mathrm{c}) \Rightarrow(\mathrm{a})$ Is just the fact that an ordering $P$ contains $\sum A^{2}$.
defn:multirealradical
Definition 1.7.2. The real radical of an ideal $\mathfrak{a}$ in $A$ is

$$
\sqrt[R]{\mathfrak{a}}:=\left\{a \in A: \exists b_{i} \in A \text { and } k \geq 0 \text { such that }\left(a^{2 k}+\sum b_{i}^{2}\right) \cap \mathfrak{a} \neq \emptyset\right\} .
$$

prop:6.2marshall
Proposition 1.7.3 (Proposition 6.2 of $[47]) . \sqrt[R]{\mathfrak{a}}$ is the intersection of all real prime ideals of $A$ containing $\mathfrak{a}$.

Proof. The inclusion $\subseteq$ is immediate because $\sqrt[R]{\mathfrak{a}}$ is real. For $\supseteq$, we use Corollary 1.6.4(a). Suppose that $a \in \mathfrak{p}$ for each real prime ideal $\mathfrak{p}$ with $\mathfrak{a} \subseteq \mathfrak{p}$. Consider $T=\sum A^{2}+\mathfrak{a}$ (the preordering in $A$ generated by $a$ ). Then $\bar{a}=0$ on $X_{T}$ so, by Corollary 1.6.4 (a), $-a^{2 k} \in T$ for some $k \geq 0$. Then $\left(a^{2 k}+\sum b_{i}^{2}\right) \cap \mathfrak{a} \neq \emptyset$ for some $b_{j}$, and $a \in \sqrt[R]{\mathfrak{a}}$.
prop:6.3marshall
Proposition 1.7.4 (Proposition 6.3 of (47]). For an ideal $\mathfrak{a}$ of a multiring $A$, the following are equivalent:
$a-\mathfrak{a}$ is real.
$b-\sqrt[R]{\mathfrak{a}}=a$.
$c-\mathfrak{a}$ is the intersection of real prime ideals.
$d-\mathfrak{a}$ is radical and every minimal prime ideal over $\mathfrak{a}$ is real.
Proof. We already have (a) $\Leftrightarrow(\mathrm{b})$, and $(\mathrm{b}) \Leftrightarrow(\mathrm{c})$ is consequence of Proposition 1.7.3. If $\mathfrak{a}$ is radical, then $\mathfrak{a}$ is the intersection of the minimal prime ideals over $\mathfrak{a}$, so $(\mathrm{d}) \Rightarrow(3)$. It remains to show that $(\mathrm{c}) \Rightarrow(\mathrm{d})$. Suppose $\mathfrak{q}$ is a minimal prime ideal over $\mathfrak{a}$ which is not real. Thus, for every real prime ideal $\mathfrak{p}$ of $A$ which $\mathfrak{a} \subseteq \mathfrak{p}$, there exists $a_{\mathfrak{p}} \in \mathfrak{p}$ such that $a_{\mathfrak{p}} \notin \mathfrak{q}$. By the compactness of $\operatorname{Sper}(A)$ in the patch topology, there exist finitely many elements $a_{1}, \ldots, a_{n}$ of $A$ such that $a_{i} \notin \mathfrak{q}$ for each $i$, and for each real prime ideal $\mathfrak{p}$ with $\mathfrak{a} \subseteq \mathfrak{p}, a_{i} \in \mathfrak{p}$ for some $i$. Let $a=a_{1} \cdot \ldots \cdot a_{n}$. Then $a \in \mathfrak{p}$ for each real prime ideal $\mathfrak{p}$ containing $\mathfrak{a}$ so, by (c), $a \in \mathfrak{a}$. This contradicts $a \notin \mathfrak{q}$.
defn:multiringreal
Definition 1.7.5. A multiring $A($ with $1 \neq 0)$ is said to be real if the ideal $\{0\}$ is real.
If $\mathfrak{a}$ is a real proper ideal of $A$, then $A / \mathfrak{a}$ is real. In particular, if $-1 \notin \sum A^{2}$, then $A / \sqrt[R]{\{0\}}$ is real.

### 1.8 Real Reduced Multirings

We assume that $A$ is a multiring with $-1 \notin \sum A^{2}$ and $T$ is a proper preordering of $A$. We use the notation of section 8.4 , where we define the multiring homomorphism $A \rightarrow Q_{2}^{X_{T}}$, given by $a \mapsto \bar{a}$, where $\bar{a}$ is defined by $\bar{a}(\sigma)=\sigma(a)$ for all $\sigma \in X_{T}$. We want to prove that the image of $A$ in $Q_{2}^{X_{T}}$ is a multiring which is strongly embedded in $Q_{2}^{X_{T}}$. Now, we will introduce some notation:
defn:multivalue
Definition 1.8.1. For $a_{1}, \ldots, a_{n} \in A$, we define the value set of $\phi=\left(\bar{a}_{1}, \ldots, \bar{a}_{n}\right)$ to be

$$
D(\phi)=D\left(\bar{a}_{1}, \ldots, \bar{a}_{n}\right)=\left\{\bar{b}: b \in \sum T a_{1}+\ldots+\sum T a_{n}\right\} .
$$

We say that $\bar{b}$ is represented by $\phi$ if $\bar{b} \in D(\phi)$.
lem:7.1marshall
Lemma 1.8.2 (Lemma 7.1 of [47).

$$
\begin{aligned}
& i-D(\bar{a})=\left\{\bar{b}^{2} \bar{a}: b \in A\right\}=\{\bar{t} \bar{a}: t \in A, \bar{t} \geq 0\}= \\
& \quad\left\{\bar{b}: \text { for each } \sigma \in X_{T} \text { either } \bar{b}(\sigma)=0 \text { or } \bar{a}(\sigma) \bar{b}(\sigma)>0\right\} . \\
& i i-D(\bar{a}, \bar{b})=\left\{\bar{c}: \text { for each } \sigma \in X_{T}, \text { either } \bar{c}(\sigma)=0 \text { or } \bar{a}(\sigma) \bar{c}(\sigma)>0 \text { or } \bar{b}(\sigma) \bar{c}(\sigma)>0\right\} \text {. } \\
& \text { iii - If } n \geq 3, D\left(\bar{a}_{1}, \ldots, \bar{a}_{n}\right)=\bigcup_{\bar{c} \in D\left(\bar{a}_{2}, \ldots, \bar{a}_{n}\right)} D\left(\bar{a}_{1}, \bar{c}\right) .
\end{aligned}
$$

iv - $D\left(\bar{a}_{1}, \ldots, \bar{a}_{n}\right)$ depends only on $\bar{a}_{1}, \ldots, \bar{a}_{n}$ (not on the particular representatives $a_{1}, \ldots, a_{n}$ ).

## Proof.

i - Is immediate from definition of $D(\bar{a})$.
ii - If $c \in \sum T a+\sum T b$, then $c^{2} \in \sum T a c+\sum T b c$. Follow this, that for any $\sigma \in X_{T}, \bar{c}(\sigma)=0$ or one of $\bar{a}(\sigma) \bar{c}(\sigma), \bar{b}(\sigma) \bar{c}(\sigma)$ is strictly positive, so $\bar{c}$ belongs to the second set. Now pick $c$ such that $\bar{c}$ belongs to the second set. Denote by $A^{\prime}$ the localization of $A$ and the multiplicative
set $S=\left\{c^{2 k} \mid k \geq 0\right\}$ and let $T^{\prime}$ be the preordering in $A^{\prime}$ defined by $T^{\prime}=\left\{t / 2^{2 k}: k \geq 0\right\}$. Let $a^{\prime}=a c, b^{\prime}=b c$. On $X_{T^{\prime}-\sum T^{\prime} a^{\prime}}, \bar{b}>0$, so by Corollary 1.6.4(b),

$$
-1 \in T^{\prime}-\sum T^{\prime} a^{\prime}-\sum A^{2} b^{\prime}
$$

Multiplying by $c^{2 m+1}$, $m$ sufficiently large, $-c^{2 m+1} \in T c-\sum T a-\sum T b$. This yields

$$
c_{1} \in\left(\sum T a+\sum T b\right) \cap\left(c^{2 m+1}+T c\right)
$$

It follows that $\bar{c}=\bar{c}_{1} \in D(\bar{a}, \bar{b})$.
iii - This follows from (ii) by induction. Note that $D(\bar{a}, \bar{c})$ depends only on $\bar{c}$, not on the particular representative of $c$.
iv - For $n=1$ and 2 , this is immediate from (i) and (ii). For $n \geq 3$, it follows by induction on $n$ using (iii).
lem:7.2marshall
Lemma 1.8.3 (Lemma 7.2 of [47]). For $a_{0}, \ldots, a_{n} \in A$, the following are equivalent:
$i$ - There exists $a_{i}^{\prime} \in A$ such that ${\overline{a^{\prime}}}_{i}=\bar{a}_{i}$ and $0 \in a_{0}^{\prime}+\ldots+a_{n}^{\prime}$.
$i i-\overline{-a}_{i} \in D\left(\bar{a}_{1}, \ldots, \bar{a}_{i-1}, \bar{a}_{i+1}, \ldots, \bar{a}_{n}\right)$ for $i=0, \ldots, n$.
Proof. (i) $\Rightarrow$ (ii) By symmetry, it is suffice to show $-\bar{a}_{0} \in D\left(\bar{a}_{1}, \ldots, \bar{a}_{n}\right)$. Since $0 \in a_{0}^{\prime}+\ldots+a_{n}^{\prime}$, $-a_{0}^{\prime} \in a_{1}^{\prime}+\ldots+a_{n}^{\prime}$, so $\bar{a}_{0}={\overline{a^{\prime}}}_{0} \in D\left({\overline{a^{\prime}}}_{1}, \ldots,{\overline{a^{\prime}}}_{n}\right)=D\left(\bar{a}_{1}, \ldots, \bar{a}_{n}\right)$, using Lemma 1.8.2(iii).
(ii) $\Rightarrow$ (i) We have $a_{i}^{\prime}$ with ${\overline{a^{\prime}}}_{i}=\bar{a}_{i}$ such that $0 \in a_{i}^{\prime}+\sum_{i \neq j} \sum T a_{j}$. Then

$$
0 \in 0+\ldots+0 \subseteq \sum_{i=0}^{n}\left(a_{i}^{\prime}+\sum_{i \neq j} \sum T a_{j}\right)=\sum_{i=0}^{n}\left(a_{i}^{\prime}+\sum T a_{i}\right)
$$

so there exist $a_{i}^{\prime \prime} \in a_{i}^{\prime}+\sum T a_{i}$ such that $0 \in a_{0}^{\prime \prime}+\ldots+a_{n}^{\prime \prime}$. Hence ${\overline{a^{\prime \prime}}}_{i}=\bar{a}_{i}$.

Denote the image of $A$ in $Q_{2}^{X_{T}}$ by $Q_{T}(A)$. Addition on $Q_{T}(A)$ is defined by $\bar{a}+\bar{b}:=\{\bar{c}: c \in a+b\}$, $\bar{a} \bar{b}:=\overline{a b},-\bar{a}:=\overline{-a}$. The zero element of $Q_{T}(A)$ is $\overline{0}$.
prop:7.3marshall
Proposition 1.8.4 (Local-Global principle, Proposition 7.3 of [47]). Let $A$ be a multiring with $-1 \notin \sum A^{2}$ and $T$ a proper preordering of $A$. Then:
$i-Q_{T}(A)$ is a multiring.
ii - $Q_{T}(A)$ is strong embedded in $Q_{2}^{X_{T}}$.

## Proof.

i - Everything is straightforward calculations except the associativity. Let $x, u, v, w, p \in A$ such that $\bar{p} \in \bar{u}+\bar{v}$ and $\bar{x} \in \bar{p}+\bar{w}$. Then $\bar{x} \in D(\bar{p}, \bar{w})$ and $\bar{p} \in D(\bar{u}, \bar{v})$, so $\bar{x} \in D(\bar{u}, \bar{v}, \bar{w})$. Also $-\bar{w} \in-\bar{x}+\bar{p}$, so $-\bar{w} \in D(-\bar{x}, \bar{p})$, i.e, $-\bar{w} \in D(-\bar{x}, \bar{u}, \bar{v})$. Also $-\bar{u} \in-\bar{p}+\bar{v}$ and $-\bar{p} \in-\bar{x}+\bar{w}$, so $-\bar{u} \in D(-\bar{p}, \bar{v})$ and $-\bar{p} \in D(-\bar{x}, \bar{w})$ i.e., $-\bar{u} \in D(-\bar{x}, \bar{v}, \bar{w})$. According to Lemma 1.8.3, this implies there exist $x^{\prime}, u^{\prime}, v^{\prime}, w^{\prime} \in A$ such that $\overline{x^{\prime}}=\bar{x}, \overline{u^{\prime}}=\bar{u}, \overline{v^{\prime}}=\bar{v}, \overline{w^{\prime}}=\bar{w}$ and $x^{\prime} \in u^{\prime}+v^{\prime}+w^{\prime}$. Pick $q \in v^{\prime}+w^{\prime}$ such that $x^{\prime} \in u^{\prime}+q$. Then $\bar{q} \in \bar{v}+\bar{w}$ and $\bar{x} \in \bar{u}+\bar{q}$.
ii - Let $a, b, c \in A$. According to Lemma 1.8.3, $\bar{c} \in \bar{a}+\bar{b}$ iff $\bar{c} \in D(\bar{a}, \bar{b}),-\bar{a} \in D(-\bar{c}, \bar{b})$ and $-\bar{b} \in D(-\bar{c}, \bar{a})$. According to Lemma 1.8 .2 (ii), this occurs iff for all $\sigma \in X_{T}, \bar{c}(\sigma) \bar{a}(\sigma)>0$ or $\bar{c}(\sigma) \bar{b}(\sigma)>0$ or $\bar{a}(\sigma) \bar{b}(\sigma)<0$ or $\bar{a}(\sigma)=\bar{b}(\sigma)=\bar{c}(\sigma)=0$, i.e., iff for all $\sigma \in X_{T}$, $\bar{c}(\sigma) \in \bar{a}(\sigma)+\bar{b}(\sigma)$.

The real spectrum of $Q_{T}(A)$ is naturally identified with $X_{T}$. Now that we know that addition is a well-defined associative operation on subsets of $Q_{T}(A)$, we have another more intrinsic description of value sets:
cor:7.4marshall
Corollary 1.8.5 (Corollary 7.4 of [47]). Let $\bar{T}=\{\bar{t}: t \in T\}=\{\bar{t}: t \in A, \bar{t} \geq 0\}$. Then:

$$
i-\bar{T} \bar{a}_{1}+. .+\bar{T} \bar{a}_{n}=\left\{\bar{b}: b \in \sum T a_{1}+\ldots+\sum T a_{n}\right\} .
$$

$i i-\overline{0} \in \bar{a}_{1}+\ldots+\bar{a}_{n} \Leftrightarrow-\bar{a}_{i} \in \sum_{j \neq i} \bar{T} \bar{a}_{j}$, for $i=0, \ldots, n \Leftrightarrow$ there exists $a_{0}^{\prime}, \ldots, a_{n}^{\prime}$ such that $0 \in a_{1}^{\prime}+\ldots a_{n}^{\prime}$ and $\bar{a}_{i}^{\prime}=\bar{a}_{i}, i=0, \ldots, n$.
Proof. (i) is direct consequence of Lemma 1.8 .2 and (ii) is direct consequence of 1.8.3

We restrict our attention now to the case where $T=\sum A^{2}$ and consider the multiring morphism $a \mapsto \bar{a}$ from $A$ into $Q_{2}^{\operatorname{Sper}(A)}$. We denote $Q_{\sum A^{2}}(A)$ by $Q_{\text {red }}(A)$ which we refer to as the real reduced multiring associated to $A$. The multirings $A$ such that the morphism $A \rightarrow Q_{r e d}(A)$ is an isomorphism are obviously of special interest.
prop:7.5marshall
Proposition 1.8.6 (Proposition 7.5 of [47]). For a multiring A with $-1 \notin \sum A^{2}$, the map $a \mapsto \bar{a}$ from $A$ onto $Q_{r e d}(A)$ is an isomorphism if and only if $A$ satisfies the following properties:
$a-a^{3}=a$.
$b-a+a b^{2}=\{a\}$.
$c-a^{2}+b^{2}$ contains a unique element.
Proof. ( $\Rightarrow$ ) By construction we have (a) and (b) (since $\bar{a}+\bar{a}=\bar{a}$ and $\bar{b}^{2}=\overline{1}$ or $\bar{b}^{2}=0$ in $Q_{\text {red }}(A)$ ). For (c), if $c \in a^{2}+b^{2}$, then $c^{2} \in\left(a^{2}+b^{2}\right)\left(a^{2}+b^{2}\right) \subseteq a^{4}+a^{2} b^{2}+a^{2} b^{2}+b^{4}=\left(a^{2}+a^{2} b^{2}\right)+\left(b^{2}+a^{2} b^{2}\right)$. Since $a^{2}+a^{2} b^{2}=\left\{a^{2}\right\}$ and $b^{2}+a^{2} b^{2}=\left\{b^{2}\right\}$, this implies $c^{2} \in a^{2}+b^{2}$. Consequently, $c^{2}=c$, i.e., the unique element of $a^{2}+b^{2}$ is necessarily a square. It follows by induction that, for any $a_{1}, \ldots, a_{n} \in A, a_{1}^{2}+\ldots+a_{n}^{2}$ contains a unique element, which is a square. In particular, $\sum A^{2}=A^{2}$.
$(\Leftrightarrow)$ Let $T=\sum A^{2}=A^{2}$. suppose that $\bar{a}=\bar{b}$. Let $c \in a^{2}+b^{2}$. Thus $-c^{2 k} \in A^{2}-\sum A^{2} a b$. Since $c^{3}=c, c^{2 k}=c^{2}$. Thus, there exists $d \in \sum A^{2} a b$ with $d \in c^{2}+A^{2}$. Hence

$$
a c \in a\left(a^{2}+b^{2}\right) \subseteq a^{3}+a b^{2}=a+a b^{2}=a,
$$

so $a c=a$. Similarly, $b c=b$ and $c d=c$. Thus, $a d=(a c) d=a(c d)=a c=a$ and, similarly, $b d=b$. Say $d \in \sum e_{i}^{2} a b$. Then $a b=a b d \in \sum e_{i}^{2} a^{2} b^{2} \subseteq A^{2}$. This implies $a b \in A^{2}$, so $a b=a^{2} b^{2}$. Thus, $a^{2}=a^{2} d \in \sum e_{i}^{2} a^{3} b=\sum e_{i}^{2} a b=\sum e_{i}^{2} a^{2} b^{2}$ and, similarly, $b^{2} \in \sum e_{i}^{2} a^{2} b^{2}$. Since $\sum e_{i}^{2} a^{2} b^{2}$ is a singleton set, this implies $a^{2}=a b=b^{2}$. Finally,

$$
a=a^{3}=a a^{2}+a b^{2}=a(a b)+a b^{2}=a^{2} b+a b^{2}=(a b) b+a b^{2}=(a b) b=b^{2} b=b^{3}=b,
$$

as required.

Definition 1.8.7. A multiring satisfying $-1 \notin \sum A^{2}$ and the equivalent conditions of Proposition 1.8.6 will be called real reduced multiring. A morphism of real reduced multirings is just a morphism of multirings. The category of real reduced multirings will be denoted by $\mathcal{M} \mathcal{R}_{\text {red }}$.

## cor:7.6marshall

Corollary 1.8.8 (Corollary 7.6 of [47]). A multiring $A$ is real reduced if and only if the following properties holds for all $a, b, c, d \in F$ :

$$
\begin{aligned}
& i-1 \neq 0 \\
& i i-a^{3}=a \\
& i i i-c \in a+a b^{2} \Rightarrow c=a \\
& i v-c \in a^{2}+b^{2} \text { and } d \in a^{2}+b^{2} \text { implies } c=d
\end{aligned}
$$

Proof. As noted above, (ii),(iii) and (iv) imply $\sum A^{2}=A^{2}$. If $-1 \in \sum A^{2}$, then $-1=a^{2}$ for some $a$, so $0 \in 1+a^{2}$. By (iii), $0=1$ and this contradicts (i). Thus $-1 \notin \sum A^{2}$. Now apply Proposition 1.8 .6 to conclude that $A$ is a real reduced multiring. The converse is immediate.

## Chapter 2

## Hyperfields, Special Groups and Quadratic Forms

There are many of abstract theories of quadratic forms. The first ones (abstract Witt rings, quaternionic structures and Cordes schemes [46]) have appeared in the late 70s, by the hands of M. Marshall and C. M. Cordes, with the following central target: analyze the existence (or not) of fields with certain properties relating to quadratic forms. In the decade of 80 's, appears the Marshall's abstract space of orderings (AOS) [48]: they are important because generalize both theory of orderings on fields and the reduced theory of quadratic forms. But only in the early 90's that arise a (finitary) first-order theory that generalizes the reduced and non-reduced theory of quadratic forms simultaneously. This theory is the special groups of F. Miraglia and M. Dickmann [28]. At that moment, the focus was to look at generalizations for the theory of quadratic forms with invertibles coefficients (fields, von Neumman rings, semi-local rings..., in general, rings with a good amount of invertibles). In the mid 90's, Marshall generalizes the abstract ordering spaces to rings, and called his new theory by "abstract real spectra" (ARS), in a first attempt to develop a theory of quadratic forms over (general) coefficients on rings. The ring-theoretic case is much more difficult to deal than the field one, the isometry is not well behaved and an algebraic counterpart of the abstract real spectra just appears in years 2000, with the real semigroups (RS) of Dickmann and Petrovich.

Following the work of professors F. Miraglia and M. Dickmann, through a fruitful and successful partnership between IME-USP and IMJ-PRG (Paris 6,7), which began in the 1990s, the three authors of this paper continue to expand the boundaries of abstract theories of quadratic forms, carrying forward the ideas of Dickmann-Miraglia's works, making the IME-USP a center for the development of such theories.

All those abstract theories constitute categories that are equivalent, or dually equivalent to full subcategories of each other. Also, each one has a particular motivation and advantage. In particular, some of them are categories of first-order theories and the corresponding language homomorphisms, thus allowing the application of model-theoretical notions and methods in this subject of algebra.

In [28], [32] and [33] are considered special groups and real semigroups. The former treats simultaneously reduced and non-reduced theories but focuses on rings with a good amount of invertible coefficients to quadratic forms. The latter has the advantage of potentially consider general coefficients of a ring, but only addresses the reduced case. Both are first-order theory, thus they allow the use of model theoretic methods.
M. Marshall in 47] introduced an approach to (reduced) theory of quadratic forms trough the concept of multiring: this seems more intuitive for an algebraist, encompassing some techniques of
ordinary commutative algebra, encodes copies of special groups and real semigroups (see [24]), but still allows the use of model-theoretic tools.

In this Chapter we study the relations between special groups, real semigroups and multivalued structures. The main results are Theorem $2.3 .4,2.3 .7,2.3 .10$ and 2.5 .4 , which characterize precisely the necessary conditions for a hyperfield/multiring come from a special group/real semigroup. Proposition 2.4.3 deals with a question posed by the authors of [37. We also got a new and interesting Example of real semigroup (2.5.15): $A /{ }_{m} T$ for $A=\mathcal{C}(X, \mathbb{R})$ and $T=A^{2} \cap \operatorname{nzd}(A)$, where $X$ is a $T_{6}$ topological space.

### 2.1 Special Groups

Definition 2.1.1 (Extension of a Relation). Let $A$ be a set and $\equiv$ a binary relation on $A \times A$. We extend $\equiv$ to a binary relation $\equiv_{n}$ on $A^{n}$, by induction on $n \geq 1$, as follows:
$i-\equiv_{1}$ is the diagonal relation $\Delta_{A} \subseteq A \times A$
$i i-\equiv_{2}=\equiv$.
iii - if $n \geq 3,\left\langle a_{1}, \ldots, a_{n}\right\rangle \equiv_{n}\left\langle b_{1}, \ldots, b_{n}\right\rangle$ if and only there are $x, y, z_{3}, \ldots, z_{n} \in A$ such that

$$
\begin{aligned}
& \left\langle a_{1}, x\right\rangle \equiv\left\langle b_{1}, y\right\rangle \\
& \left\langle a_{2}, \ldots, a_{n}\right\rangle \equiv_{n-1}\left\langle x, z_{3}, \ldots, z_{n}\right\rangle \text { and } \\
& \left\langle b_{2}, \ldots, b_{n}\right\rangle \equiv_{n-1}\left\langle y, z_{3}, \ldots, z_{n}\right\rangle
\end{aligned}
$$

Whenever clear from the context, we frequently abuse notation and indicate the aforedescribed extension $\equiv$ by the same symbol.
defn:sg
Definition 2.1.2 (Special Group, 1.2 of [28]). A special group is an tuple ( $G,-1, \equiv$ ), where $G$ is a group of exponent 2, i.e, $g^{2}=1$ for all $g \in G ;-1$ is a distinguished element of $G$, and $\equiv \subseteq G \times G \times G \times G$ is a relation (the special relation), satisfying the following axioms for all $a, b, c, d, x \in G$ :

SG $\mathbf{0} \equiv$ is an equivalence relation on $G^{2}$;
SG $1\langle a, b\rangle \equiv\langle b, a\rangle ;$
SG $2\langle a,-a\rangle \equiv\langle 1,-1\rangle$;
SG $3\langle a, b\rangle \equiv\langle c, d\rangle \Rightarrow a b=c d$;
SG $4\langle a, b\rangle \equiv\langle c, d\rangle \Rightarrow\langle a,-c\rangle \equiv\langle-b, d\rangle$;
SG $5\langle a, b\rangle \equiv\langle c, d\rangle \Rightarrow\langle g a, g b\rangle \equiv\langle g c, g d\rangle$, for all $g \in G$.
SG 6 (3-transitivity) the extension of $\equiv$ for a binary relation on $G^{3}$ (as in 2.1.1) is a transitive relation.

A group of exponent 2, with a distinguished element -1 , satisfying the axioms SG0-SG3 and SG5 is called a proto special group; a pre special group is a proto special group that also satisfies SG4. Thus a special group is a pre-special group that satisfies SG6 (or, equivalently, for each $n \geq 1, \equiv_{n}$ is an equivalence relation on $G^{n}$.)

A $n$-form (or form of dimension $n \geq 1$ ) is an $n$-tuple of elements of a pre-special group $G$. An element $b \in G$ is represented on $G$ by the form $\varphi=\left\langle a_{1}, \ldots, a_{n}\right\rangle$, in symbols $b \in D_{G}(\varphi)$, if there exists $b_{2}, \ldots, b_{n} \in G$ such that $\left\langle b, b_{2}, \ldots, b_{n}\right\rangle \equiv \varphi$.

A pre-special group (or special group) $(G,-1, \equiv)$ is:

- formally real if $-1 \notin \bigcup_{n \in \mathbb{N}} D_{G}(n\langle 1\rangle)$;
- reduced if it is formally real and, for each $a \in G, a \in D_{G}(\langle 1,1\rangle)$ iff $a=1$.

Now, some examples:
ex2.2
Example 2.1.3 (The trivial special relation, 1.9 of [28]). Let $G$ be a group of exponent 2 and take -1 as any element of $G$ different of 1. For $a, b, c, d \in G$, define $\langle a, b\rangle \equiv_{t}\langle c, d\rangle$ if and only if $a b=c d$. Then $G_{t}=\left(G, \equiv_{t},-1\right)$ is a special group ([28]). In particular $2=\{-1,1\}$ is a reduced special group.

Example 2.1.4 (Special group of a field, Theorem 1.32 of [28]). For $F$ be a field, denote
$\dot{F}=F \backslash\{0\}, \dot{F}^{2}=\left\{x^{2}: x \in \dot{F}\right\}$ and $\Sigma \dot{F}^{2}=\left\{\sum_{i \in I} x_{i}^{2}: I\right.$ is finite and $\left.x_{i} \in \dot{F}^{2}\right\} . \operatorname{Let} G(F)=\dot{F} / \dot{F}^{2}$. In the case of $F$ is be formally real, we have $\Sigma \dot{F}^{2}$ is a subgroup of $\dot{F}$, then we take $G_{r e d}(F)=$ $\dot{F} / \Sigma \dot{F}^{2}$. Note that $G(F)$ and $G_{r e d}(F)$ are groups of exponent 2. In [28] they prove that $G(F)$ and $G_{r e d}(F)$ are special groups with the special relation given by usual notion of isometry (see for instance, 43]), and $G_{r e d}(F)$ is always reduced.
defnmorph
Definition 2.1.5 (1.1 of [28]). A map $\left(G, \equiv_{G},-1\right) \xrightarrow{f}\left(H, \equiv_{H},-1\right)$ between pre-special groups is a morphism of pre-special groups or PSG-morphism if $f: G \rightarrow H$ is a homomorphism of groups, $f(-1)=-1$ and for all $a, b, c, d \in G$

$$
\langle a, b\rangle \equiv_{G}\langle c, d\rangle \Rightarrow\langle f(a), f(b)\rangle \equiv_{H}\langle f(c), f(d)\rangle
$$

A morphism of special groups or $\boldsymbol{S G}$-morphism is a $p S G$-morphism between the correspondents pre-special groups. $f$ will be an isomorphism if is bijective and $f, f^{-1}$ are PSG-morphisms.

It can be verified that a special group $G$ is formally real iff it admits some SG-morphism $f: G \rightarrow 2$.

The category of special groups (respectively reduced special groups) and theirs morphisms will be denoted by $S G$ (respectively $R S G$ ). Now, we will analyze the connections between the $S G$ and $M F$. For this, we need more results about special groups and their characterization. For this, we use the results proved in Lira's thesis [22]. Consider these axioms concerns about a group of exponent 2 with a distinguished element:

SG $7 \forall a \forall a^{\prime} \forall x \forall t \forall t^{\prime} \forall y\left[\left(a, a^{\prime}\right) \equiv(x, t) \wedge\left(t, t^{\prime}\right) \equiv(1, y)\right]$
$\Rightarrow \exists a^{\prime \prime} \exists s \exists s^{\prime}\left[\left(a, a^{\prime \prime}\right) \equiv(y, s) \wedge\left(s, s^{\prime}\right) \equiv(1, x)\right]$.
An equivalent statement for SG7 is

$$
\bigcup_{t \in D_{G}(1, y)} D_{G}(x, t)=\bigcup_{s \in D_{G}(1, x)} D_{G}(y, s)
$$

for all $x, y \in G$.
SG 8 For all forms $f_{1}, \ldots, f_{n}$ of dimension 3 and for all $a, a_{2}, a_{3}, b_{2}, b_{3} \in G$,

$$
\left\langle a, a_{2}, a_{3}\right\rangle \equiv f_{1} \equiv \ldots \equiv f_{n} \equiv\left\langle a, b_{2}, b_{3}\right\rangle \Rightarrow\left\langle a_{2}, a_{3}\right\rangle \equiv\left\langle b_{2}, b_{3}\right\rangle
$$

SG $9 \forall a \forall b \forall c \forall d[\langle a, b, a b\rangle \equiv\langle c, d, c d\rangle \Rightarrow\langle a, b, a b\rangle \equiv\langle d, c, c d\rangle]$
Proposition 2.1.6 (A. de Lima, [22]). Let $(G,-1, \equiv)$ be a pre-special group. The following are equivalent:

$$
\begin{aligned}
i-G & \models S G 6 \\
i i-G & \models S G 7 \wedge S G 8 \\
i i i-G & \models S G 9
\end{aligned}
$$

Theorem 2.1.7. Let $(G, \equiv,-1)$ be a pre-special group. The following are equivalent:
$a-\equiv$ is 3-transitive (i.e, transitive for 3-forms, and hence $G$ is a special group).
$b-\equiv$ is transitive (i.e, transitive for $n$-forms for all $n \geq 2$ ).
$c$ - For all $n \geq 2$, for all $n$-forms $\varphi, \psi$ over $G$ and all $\sigma \in S_{n}$,

$$
\varphi \equiv \psi \text { implies } \varphi \equiv \psi^{\sigma} .
$$

$d$ - For all $n \geq 2$, for all $n$-forms $\varphi, \psi$ over $G$,

$$
\varphi \equiv \psi \text { iff } \varphi \approx \psi
$$

$e$ - For all 3 -forms $\varphi$ and all $b_{1}, b_{2}, b_{3} \in G$,

$$
\varphi \equiv\left\langle b_{1}, b_{2}, b_{3}\right\rangle \text { implies } \varphi \equiv\left\langle b_{2}, b_{1}, b_{3}\right\rangle
$$

1.24 chico

Corollary 2.1.8. Let $(G, \equiv,-1)$ be a pre-special group, $\varphi$ and $\psi$ be forms over $G$ and $a, b, x, y \in G$. The following are equivalent:
$a-G$ is a special group.
$b$ - For all forms $\varphi, \psi$ over $G$ and all $a, b, x, y \in G$

$$
\varphi \equiv\langle a, b\rangle \oplus \psi \text { and }\langle a, b\rangle \equiv\langle x, y\rangle \Rightarrow \varphi \equiv\langle x, y\rangle \oplus \psi .
$$

$c$ - For all 3 -forms $\varphi, \psi$ over $G$ and all $a, b, c, x, y \in G$

$$
\varphi \equiv\langle a, b, c\rangle \text { and }\langle a, b\rangle \equiv\langle x, y\rangle \Rightarrow \varphi \equiv\langle x, y, c\rangle .
$$

2.3chico

Definition 2.1.9 (2.3 of [28]). Let $G$ be a special group and let $\Delta \subseteq G$ be a subgroup. We say that $\Delta$ is saturated if for all $a \in G$,

$$
\begin{equation*}
a \in \Delta \Rightarrow D_{G}(1, a) \subseteq \Delta . \tag{sat}
\end{equation*}
$$

Note that if, in addition, $-1 \in \Delta$, then $\Delta=G$. Thus we will reserve the noun saturated for those subgroups satisfying [sat] such that $-1 \notin \Delta$, while $G$ will be called the improper saturated subgroup of itself.

Lemma 2.1.10 (2.4 of [28]). Let $G$ be a special group and $\Delta$ a subgroup of $G$.
a - The intersection of any family of saturated subgroups is saturated. The union of an upward directed family of saturated subgroup is saturated.
$b$ - The following are equivalent:
$i-\Delta$ is saturated.
ii - For any Pfister forms $\varphi, \psi$ over $\Delta$ and any $b, c \in \Delta$

$$
D_{G}(\varphi), D_{G}(\psi) \subseteq \Delta \Rightarrow D_{G}(b \varphi \oplus c \psi) \subseteq \Delta
$$

iii - For any Pfister form $\varphi$ over $\Delta, D_{G}(\varphi) \subseteq \Delta$.

### 2.2 Special Hyperfields

sg.to.mf
Proposition 2.2.1. Let $(G, \equiv,-1)$ be a special group and $M(G):=G \cup\{0\}$ where $0:=\{G\}$. Then ( $M(G),+,-, \cdot, 0,1)$ with operations

- $a \cdot b=\left\{\begin{array}{l}0 \text { if } a=0 \text { or } b=0 \\ a \cdot b \text { otherwise }\end{array}\right.$
- $-(a)=(-1) \cdot a$
- $a+b=\left\{\begin{array}{l}\{b\} \text { if } a=0 \\ \{a\} \text { if } b=0 \\ M(G) \text { if } a=-b, \text { and } a \neq 0 \\ D_{G}(a, b) \text { otherwise }\end{array}\right.$
is a hyperfield.
Proof. Firstly, note that + is well-defined. Then, we verify the conditions of Definition 1.2.7.
i - For this, we check the conditions of definition 1.2.1.
a $-d \in a+0=\{a\}$ imply $d=a$, and then $a \in d+(-0)$ and $0 \in(-a)+d$. Let $a=-b$ and $d \in a+(-a)=M(G)$. If $d=0$, then $a \in d+(-(-a))=0+a$ and $-a \in(-a)+0$. If $d \neq 0$, then $a \in D_{G}(d, a)$ and $-a \in D_{G}(-a, d)$ so $a \in d+(-(-a))=d+a$ and $-a \in(-a)+d$. Finally, let $a, b \neq 0$ with $a \neq-b$, and $d \in a+b$. Then there exist $g \in M(G) \backslash\{0\}$ such that $\langle d, g\rangle \equiv\langle a, b\rangle$. By SG4, $\langle d,-a\rangle \equiv\langle-g, b\rangle$ (and $\langle b,-g\rangle \equiv\langle-a, d\rangle$ by SG1). So $a \in d+(-b)$ and $b \in(-a)+d$.
$\mathrm{b}-(y \in x+0) \Leftrightarrow(x=y)$ is an immediate consequence of the definition of sum.
c $-a+0=0+a$ and $a+(-a)=M(G)=(-a)+a$. Let $a, b \in M(G), a, b \neq 0$ and $a \neq-b$. Since $D_{G}(a, b)=D_{G}(b, a)$, we have $a+b=b+a$ proving the commutativity. Observe that if $a, b \neq 0$ with $a \neq-b$, then $0 \notin a+b$.
d - Now we prove the associativity. Let $a=0$ (the cases $b=0$ and $c=0$ are analogous).
Then $0+(b+c)=\{0+g: g \in b+c\}=b+c$ and $(0+b)+c=(\{b\})+c=b+c$.
Now, let $a, b, c \neq 0$ with $a=-c$.

$$
(a+b)+(-a)=\bigcup\{g+(-a): g \in a+b\}=M(G)(\mathrm{I})
$$

because $a \in a+b$, and

$$
a+(b+(-a))=\bigcup\{a+h: h \in b+(-a)\}=M(G)(\mathrm{II})
$$

because $-a \in b+(-a)$. So (I) $=$ (II) and $(a+b)+(-a)=a+(b+(-a))$. For the case $a, b, c \neq 0, a=-b$ (the cases $b \neq-c$ is analogous) we have

$$
\begin{equation*}
(a+(-a))+c=\bigcup\{g+c: g \in M(G)\}=M(G) \tag{III}
\end{equation*}
$$

and

$$
a+((-a)+c)=\bigcup\{a+h: h \in(-a)+c\}=M(G)(\mathrm{IV})
$$

because $-a \in(-a)+c$. So (III) $=$ (IV) and $(a+(-a))+c=a+((-a)+c)$. Finally, let $a, b, c \neq 0, a \neq-b, b \neq-c$ and $a \neq-c$.

$$
(a+b)+c=c+(a+b)=\bigcup\{c+g: g \in a+b\}=\bigcup_{g \in D_{G}(a, b)} D_{G}(c, g)(\mathrm{V})
$$

and

$$
a+(b+c)=\bigcup\{h+a: h \in b+c\}=\bigcup_{h \in D_{G}(b, c)} D_{G}(h, a)(\mathrm{VI})
$$

By SG7 (applying SG5) we have $(\mathrm{V})=(\mathrm{VI})$. Then $(a+b)+c=a+(b+c)$ for all $a, b, c \in M(G)$.
ii - We conclude that $(M(G), \cdot, 1)$ is a commutative monoid as consequence of $(G, \cdot, 1)$ being an abelian group and the extended definition of • to $M(G)$. Beyond this, we have that every nonzero element of $M(G)$ has an inverse.
iii - $a \cdot 0=0$ for all $a \in M(G)$ is a consequence of the extended definition of multiplication to $M(G)$.
iv - If $a=0$ or $a \neq-b$, then $(d \in a+b) \Rightarrow \forall g(g d \in g a+g b)$ is direct consequence of the definition of sum. Next this, let $a, b \neq 0$ with $a \neq-b$ and $d \in a+b$. By SG5 $g d \in g a+b g$. Thus we have $g(a+b) \subseteq g a+g b$ for all $a, b, g \in M(G)$.

Corollary 2.2.2. The correspondence $G \mapsto M(G)$ defines a full and faithful functor

$$
M: S G \rightarrow M F
$$

Proof. Let $f: G \rightarrow H$ be a SG-morphism. We extend $f$ to $M(f): M(G) \rightarrow M(H)$ by $M(f) 1_{G}=f$ and $M(f)(0)=0$. By the definition of SG-morphism we have $M(f)(1)=1, M(f)(-a)=-a$ and $M(f)(a b)=M(f)(a) M(f)(b)$. Since $d \in D_{G}(a, b)$ implies $f(d) \in D_{H}(f(a), f(b))$ we have

$$
d \in a+b \text { imply } M(f)(d) \in M(f)(a)+M(f)(b) \text { for all } a, b \in M(G),
$$

so $M(f)$ is a multiring morphism. Now, let $f: G \rightarrow H$ and $g: H \rightarrow K$ be SG-morphisms. Since $M(f \circ g) 1_{G}=f \circ g=M(f) 1_{G} \circ M(g) 1_{G}$ and $M(f \circ g)(0)=0=M(f) \circ M(g)(0)$, we have $M(f \circ g)=M(f) \circ M(g)$. Then $M: S G \rightarrow M F$ is a functor. This functor is faithful, because if $G$ and $H$ are special groups and $f, g: G \rightarrow H$ are SG-morphisms such that
$M(f), M(g): M(G) \rightarrow M(H)$ are equal, then

$$
\left.M(f)\right|_{M(G) \backslash\{0\}}=\left.M(g)\right|_{M(G) \backslash\{0\}}
$$

and therefore $f=g$, since $M(G) \backslash\{0\}=G$.
Proposition 2.2.3. Let $G$ be an $S G$ and $M(G)$ as above. Then:
$i-a^{2}=1$ for all $a \in M(G) \backslash\{0\} ;$
ii - $1 \in 1+a$ for all $a \in M(G)$;
iii $-1+a$ is closed by multiplication for all $a \in M(G)$;
$i v$ - If there exists $x, y, z \in \dot{M}(G)$ such that

$$
\left\{\begin{array} { l } 
{ a x = c y } \\
{ a = x z } \\
{ d = y z }
\end{array} \quad \text { and } \left\{\begin{array}{l}
a \in c+y \\
b \in x+z \\
c \in y+z
\end{array}\right.\right.
$$

then there exists $t, v, w \in \dot{M}(G)$ such that

$$
\left\{\begin{array} { l } 
{ b t = c v } \\
{ b = t w } \\
{ c = v w }
\end{array} \quad \text { and } \left\{\begin{array}{l}
b \in c+v \\
a \in t+w \\
d \in v+w
\end{array}\right.\right.
$$

Proof.
i - Is just the fact of $G$ be a group of exponent 2 .
ii - Follow immediately.
iii - If $a=0$ or $a=-1$ it is trivial. If $a \neq 0,-1$, given $x, y \in 1+a=D_{G}(1, a)$, we have $\langle x, x a\rangle \equiv\langle 1, a\rangle$ and $\langle y, y a\rangle \equiv\langle 1, a\rangle$. Multiplying the first equality by 1 , we get

$$
\langle x y, x y a\rangle \equiv\langle y, y a\rangle \equiv\langle 1, a\rangle
$$

and then $x y \in D_{G}(1, a)=1+a \equiv_{G}$.
iv - Follow from 3-transitivity.
defn:special.mf
Definition 2.2.4. A hyperfield $F$ satisfying the properties $i$-iv of Proposition 2.2.3 will be called $a$ special hyperfield (SMF). Note that, if $G$ is a $S G$, then $M(G)$ is a $S M F$.
mf.to.special
Theorem 2.2.5. If $F$ is a special hyperfield the $(F \backslash\{0\}, \equiv,-1)$ is a special group where

$$
\langle a, b\rangle \equiv\langle c, d\rangle \text { iff } a b=c d \text { and } a \in c+d
$$

Proof. By (i), we have that $(F \backslash\{0\}, 1)$ is a group of exponent 2. Now, we check each axiom of Definition 2.1.2:

SG0 - By (ii) $1 \in 1+a b$, so $a b \in 1+a b$ and $a \in b+a$. Since $a b=a b$, then $\langle a, b\rangle \equiv\langle a, b\rangle$, i.e, the relation $\equiv$ is reflexive. If $\langle a, b\rangle \equiv\langle c, d\rangle$, then $a b=c d$ and $a \in c+d$. Then $a b \in c b+d b$, so by $a b=c d$, we have $c d \in a d+d b$ and then $c \in a+b$. So $\langle c, d\rangle \equiv\langle a, b\rangle$ and $\equiv$ is symmetric. Finally, suppose that $\langle a, b\rangle \equiv\langle c, d\rangle$ and $\langle c, d\rangle \equiv\langle e, f\rangle$. First, $a b=c d$ and $c d=e f$ implies $a b=e f$. Second, in order to show that $a \in e+f$, note that $a \in c+d \Rightarrow a c \in 1+c d=1+e f$ and $c \in e+f \Rightarrow c e \in 1+e f$; then by (iii), we have $a e \in 1+e f$ and so $a \in e+f$. Therefore $\langle a, b\rangle \equiv\langle e, f\rangle$.

SG1 - As $F$ is a hyperfield, $a b=b a$. By (ii), $1 \in 1+a b$, then $a b \in 1+b a$ and $b \in a+b$. Therefore $\langle a, b\rangle \equiv\langle b, a\rangle$.

SG2 - Since $1 \in 1-a$, we have $a \in 1-1$. Therefore $\langle a,-a\rangle \equiv\langle 1,-1\rangle$.
SG3 - Follow by definition.
SG4 $-\langle a, b\rangle \equiv\langle c, d\rangle \Rightarrow a b=c d$ and $a \in c+d$.

$$
\begin{gather*}
a b=c d \Rightarrow-a b b c=-b c c d \Rightarrow-a c=-b d  \tag{2;1}\\
a \in c+d \Rightarrow a d \in 1+c d=1+a b \Rightarrow d \in a+b \Rightarrow a \in-b+d \tag{2.2}
\end{gather*}
$$

so by 2.1 and 2.2 follow that $\langle a,-c\rangle \equiv\langle-b, d\rangle$.
SG5 - $\langle a, b\rangle \equiv\langle c, d\rangle \Rightarrow a b=c d$ and $a \in c+d \stackrel{I}{\Rightarrow}(g a)(g b)=(g c)(g d)$ and $g a \in g c+g d \Rightarrow$ $\langle g a, g b\rangle \equiv\langle g c, g d\rangle$.

SG6 - We use the equivalences in Theorem 2.1.8. $\langle a, b, a b\rangle \equiv\langle c, d, c d\rangle \Rightarrow$ there exists $x, y, t \in$ $F \backslash\{0\}$ such that

$$
\left\{\begin{array} { l } 
{ \langle a , x \rangle \equiv \langle c , y \rangle } \\
{ \langle b , a b \rangle \equiv \langle x , z \rangle } \\
{ \langle d , c d \rangle \equiv \langle y , z \rangle }
\end{array} \Rightarrow \left\{\begin{array}{l}
a x=c y \text { and } a \in c+y \\
a=x z \text { and } b \in x+z \\
c=y z \text { and } d \in y+z
\end{array}\right.\right.
$$

then by (v) there exists $t, v, w \in F \backslash\{0\}$ such that

$$
\left\{\begin{array} { l } 
{ b t = c v \text { and } b \in c + v } \\
{ b = t w \text { and } a \in t + w } \\
{ d = v w \text { and } d \in v + w }
\end{array} \Rightarrow \left\{\begin{array}{l}
\langle b, t\rangle \equiv\langle c, v\rangle \\
\langle a, a b\rangle \equiv\langle t, w\rangle \\
\langle d, c d\rangle \equiv\langle v, w\rangle
\end{array}\right.\right.
$$

this implies $\langle b, a, a b\rangle \equiv\langle c, d, c d\rangle$.

Corollary 2.2.6. There is a functor $S: S M F \rightarrow S G$.
Proof. In the objects of $S M F$, we define $S(F)=F \backslash\{0\}$ since the special group as stated in Theorem 2.2.5. Now, let $\sigma: F \rightarrow K$ be a SMF-morphism. Define $S(\sigma)=\left.\sigma\right|_{F \backslash\{0\}}$. We have that $S(\sigma)$ is a group homomorphism with $S(\sigma)(-1)=-1$. If $a, b \neq 0$ and $c \in a+b, c \neq 0$, then there exists $d \in F \backslash\{0\}$ such that $\langle a, b\rangle \equiv_{S(F)}\langle c, d\rangle$, and as $c \in a+b \rightarrow \sigma(c) \in \sigma(a)+\sigma(b)$, we have $\langle\sigma(a), \sigma(b)\rangle \equiv_{S(K)}\langle\sigma(c), \sigma(d)\rangle$. Therefore:

$$
(c \in a+b \rightarrow \sigma(c) \in \sigma(a)+\sigma(b)) \Rightarrow\left(c \in D_{S(F)}(a, b) \rightarrow \sigma(c) \in D_{S(K)}(\sigma(a), \sigma(b))\right)
$$

And $S(\sigma)$ is a SG-morphism. Applying the same argument, we proof that $S(\sigma \tau)=S(\sigma) S(\tau)$. Hence, $S$ is a morphism.
teo:sgsmfequiv
Theorem 2.2.7. There exist an equivalence of categories between $S G$ and $S M F$.
Proof. By the Corollaries 2.2 .2 and 2.2 .6 , we have functors $M: S G \rightarrow S M F$ and $S: S M F \rightarrow S G$. We will proof that $M \circ S \cong I d_{S M F}$ and $S \circ M \cong I d_{S G}$.
i - $M \circ S \cong I d_{S M F}$. Let $F$ be a SMF. How $S(F)=F \backslash\{0\}$ and $M(S(F))=S(F) \cup\{0\}$, we have $M(S(F))=F$. Next, let $\sigma: F \rightarrow K$ be a SMF-morphism. We have that $S(\sigma)=\left.\sigma\right|_{F \backslash\{0\}}$ and $M(S(\sigma))$ is defined with the extension $S(\sigma)(0)=0$. Therefore $M(S(\sigma))=\sigma$ and $M \circ S \cong$ $I d_{S M F}$.
ii - $S \circ M \cong I d_{S G}$. Let $G$ be a SG. Again, $M(G)=G \cup\{0\}$ and $S(M(G))=M(G) \backslash\{0\}$. Hence $S(M(G))=G$. Next, let $f: G \rightarrow H$ be a SG-morphism. How $M(f)$ is defined with the extension $f(0)=0$ and $S(M(f))=\left.M(f)\right|_{M(G) \backslash\{0\}}$, we have that $S(M(f))=f$ and $S \circ M \cong I d_{S G}$, finalizing the proof.
psgpsmfhell
Theorem 2.2.8. Let $G$ be a pre-special group and consider $(M(G),+,-, 0,1)$, with operations defined by

- $a \cdot b=\left\{\begin{array}{l}0 \text { if } a=0 \text { or } b=0 \\ a \cdot b \text { otherwise }\end{array}\right.$

$$
\bullet a+b=\left\{\begin{array}{l}
\{b\} \text { if } a=0 \\
\{a\} \text { if } b=0 \\
M(G) \text { if } a=-b, \text { and } a \neq 0 \\
D_{G}(a, b) \text { otherwise }
\end{array}\right.
$$

- $-(a)=(-1) \cdot a$

Then $M(G)$ is a pre-special multifield. Conversely, if $F$ is a pre-special multifield then $\left(\dot{F}, \equiv_{F}\right.$ , -1 ) is a pre-special group, where

$$
\langle a, b\rangle \equiv_{F}\langle c, d\rangle \text { iff } a b=c d \text { and } a \in c+d
$$

To prove it we will need a result from [28].
1.21 chico

Lemma 2.2.9 (Lemma 1.21 of $[28])$. Let $(G, \equiv,-1)$ be a pre-special group. Let $a, b, c, x, y$ be elements of $G$ and $\varphi, \psi$ be forms over $G$. Assume that $\langle a, b\rangle \equiv\langle x, y\rangle$. Then

$$
i-\text { If } \varphi \equiv\langle a, b\rangle \text { then } \varphi \equiv\langle x, y\rangle
$$

ii - For all $\sigma \in S_{3},\langle a, b, c\rangle \equiv\langle x, y, c\rangle^{\sigma}$, where

$$
\langle x, y, c\rangle^{\sigma}:=\left\langle\sigma\left(e_{1}\right), \sigma\left(e_{2}\right), \sigma\left(e_{3}\right)\right\rangle \text { with } e_{1}=x, e_{2}=y, e_{3}=c
$$

Proof of Theorem 2.2.8. Let $F$ be a pre-special hyperfield. The $\operatorname{argument}$ to proof that $(\dot{F}, \cdot, 1, \equiv)$ is a pre-special group is the same of the proof of Theorem 3.18 in [24].

Now let $(G, \cdot, 1, \equiv)$ be a pre-special group an $M(G)$ as above. Firstly, note that by SG2 and the fact that $\equiv$ is an equivalence relation we have $a-a=1-1=M(G)$ for all $a \in G$. Moreover if $x, y \in 1+a=D_{G}(1, a)$ with $a \neq-1$, we have $\langle x, x a\rangle \equiv\langle 1, a\rangle$ and $\langle y, y a\rangle \equiv\langle 1, a\rangle$. Using SG5 (and the fact that $\equiv$ is an equivalence relation) we get that $\langle 1, a\rangle \equiv\langle x, x a\rangle$ imply

$$
\langle(x y) 1,(x y) a\rangle \equiv\langle(x y) x,(x y) x a\rangle \equiv\langle y, y a\rangle \equiv\langle 1, a\rangle
$$

proving that $x y \in D_{G}(1, a)=1+a$.
Therefore, once we verify the conditions of Definition 1.2 .7 we get that $(M(G),+, \cdot, 0,1)$ is a pre-special hyperfield.

The verification of the conditions in Definition 1.2.7 is quite straightforward except perhaps by associativity, which we will prove here. We want to show that for all $a, b, c \in M(G)$,

$$
(a+b)+c=a+(b+c)
$$

If $0 \in\{a, b, c\}$ we are done. Now let $0 \notin\{a, b, c\}$. We prove that

$$
a+(b+c)=D_{G}(a, b, c) .
$$

In fact, if $x \in a+(b+c)$ then $x \in a+y$ for some $y \in b+c$. Then we have $v, w \in G$ with

$$
\langle x, v\rangle \equiv\langle a, y\rangle \text { and }\langle y, w\rangle \equiv\langle b, c\rangle .
$$

These isometries imply that $\langle x, v, w\rangle \equiv\langle a, b, c\rangle$ and then $x \in D_{G}(a, b, c)$. Conversely, let $x \in$ $D_{G}(a, b, c)$. Then

$$
\left\langle x, z_{2}, z_{3}\right\rangle \equiv\langle a, b, c\rangle
$$

for some $z_{2}, z_{3} \in G$, and hence, there are $t_{1}, t_{2}, t_{3} \in G$ with

$$
\left\langle x, t_{1}\right\rangle \equiv\left\langle a, t_{2}\right\rangle,\left\langle z_{2}, z_{3}\right\rangle \equiv\left\langle t_{1}, t_{3}\right\rangle \text { and }\langle b, c\rangle \equiv\left\langle t_{2}, t_{3}\right\rangle .
$$

Therefore $x \in a+t_{2}$ with $t_{2} \in b+c$, so $x \in a+(b+c)$. In particular, if $\langle a, b, c\rangle \equiv\langle x, y, z\rangle$, then

$$
a+(b+c)=x+(y+z) .
$$

Since $a, b \in\langle a, b\rangle$ and $\langle a, b\rangle \equiv\langle b, c\rangle$, using Lemma 2.2.9 we have

$$
\langle a, b, c\rangle \equiv\langle c, a, b\rangle \Rightarrow a+(b+c)=c+(a+b)=(a+b)+c .
$$

Then, $(M(G),+,-, \cdot, 0,1)$ is a pre-special hyperfield.

### 2.3 A Special Group associated to domains via Marshall quotient

Let $F$ be a field. There is an almost canonical way to associate a special group to $F$ (described in Example 2.1.4): consider $G_{F}:=\dot{F} / \dot{F}^{2}$ with the isometry given by the usual isometry provide by the algebraic theory of quadratic forms. As we have already seen, $G_{F}$ is the multiplicative group of units of a special hyperfield, and in this sense,

$$
M_{F}=G_{F} \cup\{0\} \cong F /{ }_{m} \dot{F}^{2}
$$

In other words, we put in correspondence special groups and special hyperfields just adding (or erasing) a zero element.

One of the main purposes of this work is extend the above situation, $M_{A} \cong A /{ }_{m} T$, where $A$ is a commutative ring with unit and $M_{A}$ is a formally real semigroup. This section deals with the case where $A$ is a domain, i.e, rings without zero divisors. Of course, we fatally need to impose some conditions to our structures:

Definition 2.3.1. An hyperbolic multiring is a multiring $R$ such that $1-1=R$.

Note that if $R$ is hyperbolic and $a \in R^{\times}$, then $R=a-a$. For a ring $R$ (i.e, the sum is univalorated), $R$ never is hyperbolic, since $1-1=\{0\}$. However, this is not a problem, since the inclusion functor $\operatorname{Ring}_{2} \hookrightarrow M \operatorname{Ring}_{2}$ is not the most natural to be considered in the quadratic forms context. Considering the special group of a field $G(F)=\dot{F} / \dot{F}^{2}$ and its special hyperfield associated, $M(G(F))=G(F) \cup\{0\}$, we get that $M(G(F))$ is hyperbolic. Hence, the desired functor to keep in mind is $M \circ G:$ Fields $_{2} \rightarrow S M F$.

Let $R$ be a ring without zero divisors. The main goal of this section is to describe conditions for a subset $T \subseteq R \backslash\{0\}$ of $R$ in such a way that $R / m T$ is a special hyperfield and therefore, (essentially) a special group. Of course, here is an abuse of notation: when we say that " $R /{ }_{m} T$ is a special group" we mean that "the induced structure in $(R / m T) \backslash\{0\}$ provides a special group strucuture".

We we seek for inspiration in the analogous conditions for the field case (see for instance, Definition 1.28 of [28], and in particular, the "completing squares" Lemma 1.29). After months of hard work, we obtained the following Definition:

Definition 2.3.2. A Dickmann-Miraglia multiring (or DM-multiring for short) ${ }^{1}$ is a pair $(R, T)$ such that $R$ is a multiring, $T \subseteq R$ is a multiplicative subset of $R \backslash\{0\}$, and $(R, T)$ satisfy the following properties:

DM0 $R / m T$ is hyperbolic.
DM1 If $\bar{a} \neq 0$ in $R / m T$, then $\bar{a}^{2}=\overline{1}$ in $R / m T$. In other words, for all $a \in R \backslash\{0\}$, there are $r, s \in T$ such that $a r=s$.

DM2 For all $a \in R$, $(\overline{1}-\bar{a})(\overline{1}-\bar{a}) \subseteq(\overline{1}-\bar{a})$ in $R / m T$.
DM3 For all $a, b, x, y, z \in R \backslash\{0\}$, if

$$
\left\{\begin{array}{l}
\bar{a} \in \bar{x}+\bar{b} \\
\bar{b} \in \bar{y}+\bar{z}
\end{array} \quad \text { in } R /{ }_{m} T,\right.
$$

then exist $\bar{v} \in \bar{x}+\bar{z}$ such that $\bar{a} \in \bar{y}+\bar{v}$ and $\overline{v b} \in \overline{x y}+\overline{a z}$ in $R / m T$.
If $R$ is a ring, we just say that $(R, T)$ is a DM-ring, or $R$ is a DM-ring. A Dickmann-Miraglia hyperfield (or DM-hyperfield) $F$ is a hyperfield such that ( $F,\{1\}$ ) is a DM-multiring (satisfy DMODM3). In other words, $F$ is a DM-hyperfield if $F$ is hyperbolic and for all $a, b, v, x, y, z \in F^{*}$,

$$
\begin{aligned}
& i-a^{2}=1 \\
& i i-(1-a)(1-a) \subseteq(1-a) \\
& \text { iii - If }\left\{\begin{array}{l}
a \in x+b \\
b \in y+z
\end{array} \text { then exist } v \in x+z \text { such that } a \in y+v \text { and } v b \in x y+a z .\right.
\end{aligned}
$$

Remark 2.3.3. These Axioms above deserves some explanation:
$i$ - Since $R$ is a domain and $0 \notin T, \bar{a}=\overline{0}$ in $R /{ }_{m} T$ iff $a=0$.
ii - DM1 entails that $R / m T$ is a hyperfield.

[^4]iii - In DM2, the expression $(1-a)(1-a)$ means multiplication of sets, i.e,
$$
(1-a)(1-a):=\{x \cdot y: x, y \in 1-a\} .
$$
iv - Looking at the expression in DM3, from
\[

\left\{$$
\begin{array}{l}
\bar{v} \in \bar{x}+\bar{z} \\
\bar{b} \in \bar{y}+\bar{z} \quad \text { in } R / m T, \\
\bar{a} \in \bar{x}+\bar{b}
\end{array}
$$\right.
\]

and the properties of multiring, we obtain

$$
\overline{v b} \in \overline{x y}+\left(\overline{x z}+\overline{y z}+\bar{z}^{2}\right) \supseteq \overline{x y}+\bar{z}(\bar{x}+\bar{y}+\bar{z}) \text { in } R / m T
$$

and

$$
\bar{a} \in \bar{x}+\bar{b} \subseteq \bar{x}+\bar{y}+\bar{z} \text { in } R / m \text {. }
$$

Hence, we can interpret the condition $\overline{v b} \in \overline{x y}+\overline{a z}$ in $R /{ }_{m} T$ as a way of "controlling" the product $\overline{v b}$ to "not escape so much" under the set $\bar{x}+\bar{y}+\bar{z}$. In the field case (when we can "change" $\in$ by $=$ ), under the Marshall's quotient the condition M3 is not necessary (see Theorem 1.32 of [28]).
$v$ - In DM3, if $0 \in\{a, b, x, y, z\}$ the axiom is trivially valid.
Theorem 2.3.4. Let $(R, T)$ be a DM-multiring and denote $\operatorname{Sm}(R, T)=\left(R /{ }_{m} T\right)$. Then $\operatorname{Sm}(R)$ is a special hyperfield (thus $\operatorname{Sm}(R, T)^{\times}$is a special group).

Remember that a special hyperfield is a hyperfield $F$ satisfying:
SMF1 $a^{2}=1$ for all $a \in \dot{F}$;
SMF2 $1 \in 1+a$ for all $a \in F$;
SMF3 $1+a$ is closed by multiplication for all $a \in \dot{F}$;
SMF4 For all $a, b, c \in \dot{F}$,

$$
\text { If } \exists p \in \dot{F} \text { such that }\left\{\begin{array} { l l } 
{ a } & { \in c + c p } \\
{ b } & { \in p + a p } \\
{ d } & { \in p + c p . }
\end{array} \text { then } \exists l \in \dot { F } \text { such that } \left\{\begin{array}{ll}
a & \in d+d l \\
b & \in l+a l \\
c & \in l+d l .
\end{array}\right.\right.
$$

Proof of Theorem 2.3.4. The properties [SMF1]-[SMF3] are imediately consequence of the axioms of sum in a multiring and [M0]-[M2] in the Definition of DM-multirings. Then, we shall prove [SMF 4]: we will rewrite de argument of Theorem 1.32 in [28]. In order to do this, we use the language of special groups. If we prove that $R /{ }_{m} T$ is a special group, then we prove that it is a special hyperfield (since [SMF 4] is precisely the translation of the axiom [SG9] for special groups to the language of hyperfields).

Here, the special relation in $R / m T$ is defined by the rule

$$
\langle\bar{a}, \bar{b}\rangle \equiv\langle\bar{c}, \bar{d}\rangle \Leftrightarrow[\overline{a b}=\overline{c d} \text { and } \bar{a} \in \bar{c}+\bar{d}](\text { in } R / m T) .
$$

Translating this to a condition with coefficients in $R$, we have

$$
\langle\bar{a}, \bar{b}\rangle \equiv\langle\bar{c}, \bar{d}\rangle \Leftrightarrow[a b v=c d w \text { and } a r \in c s+d t] \text { for some } r, s, t, v, w \in R .
$$

Using [SMF1]-[SMF3] and the multirings properties we obtain the validity of [SG0-SG5] (for more details, see Theorem 3.18 of [24]).

Hence by 2.1.6 we only need to deal with [SG9] (see condition (5) in Theorem 1.23 of [28]), and it is enough to show that

$$
\langle\bar{a}, \bar{b}, \bar{c}\rangle \equiv\langle\bar{x}, \bar{y}, \bar{z}\rangle \text { implies }\langle\bar{a}, \bar{b}, \bar{c}\rangle \equiv\langle\bar{y}, \bar{x}, \bar{z}\rangle \text {. }
$$

Suppose $\langle\bar{a}, \bar{b}, \bar{c}\rangle \equiv\langle\bar{x}, \bar{y}, \bar{z}\rangle$. Then, there exist $\alpha, \beta, \gamma$ such that

$$
\begin{equation*}
\langle\bar{a}, \bar{\alpha}\rangle \equiv\langle\bar{x}, \bar{\beta}\rangle,\langle\bar{b}, \bar{c}\rangle \equiv\langle\bar{\alpha}, \bar{\gamma}\rangle \text { and }\langle\bar{y}, \bar{z}\rangle \equiv\langle\bar{\beta}, \bar{\gamma}\rangle . \tag{2.3}
\end{equation*}
$$

Then, there exists $p_{a}, q_{a}, r_{a}, p_{\beta}, q_{\beta}, r_{\beta} \in T$ such that

$$
\begin{align*}
& a p_{a} \in x q_{a}+\beta r_{a} .  \tag{}\\
& \beta p_{\beta} \in y q_{\beta}+z r_{\beta} . \tag{2,2.5}
\end{align*}
$$

Therefore $\bar{a} \in \bar{x}+\bar{b}$ and $\bar{b} \in \bar{y}+\bar{z}$. Applying [DM3], exists

$$
\begin{equation*}
\bar{v} \in \bar{x}+\bar{z}, \tag{}
\end{equation*}
$$

such that

$$
\begin{equation*}
\bar{a} \in \bar{y}+\bar{v} \tag{2.7}
\end{equation*}
$$

We discuss two cases.
Case I: $v=0$. Then, from equation 2.7, we have $\bar{a}=\bar{y}$. Consequently, the third isometry in equation 2.3 can be written as $\langle\bar{a}, \bar{z}\rangle \equiv\langle\bar{\beta}, \bar{\gamma}\rangle$. This isometry, the first one in equation 2.3 and [SG4] yield

$$
\langle\bar{x},-\bar{\alpha}\rangle \equiv\langle\bar{a},-\bar{\beta}\rangle \equiv\langle-\bar{z}, \bar{\gamma}\rangle,
$$

and so $\langle\bar{x},-\bar{\alpha}\rangle \equiv\langle-\bar{z}, \bar{\gamma}\rangle$. Another application of [SG4] yields $\langle\bar{x}, \bar{z}\rangle \equiv\langle\bar{\alpha}, \bar{\gamma}\rangle$, which together with the second isometry in equation 2.3 gives $\langle\bar{x}, \bar{z}\rangle \equiv\langle\bar{b}, \bar{c}\rangle$. Then, we have

$$
\langle\bar{a}, \bar{x}\rangle \equiv\langle\bar{a}, \bar{x}\rangle,\langle\bar{b}, \bar{c}\rangle \equiv\langle\bar{x}, \bar{z}\rangle, \text { and }\langle\bar{x}, \bar{z}\rangle \equiv\langle\bar{x}, \bar{z}\rangle,
$$

which shows that $\langle\bar{a}, \bar{b}, \bar{c}\rangle \equiv\langle\bar{a}, \bar{x}, \bar{z}\rangle$, as required.
Case II: $v \neq 0$. Equation 2.7 implies $\bar{a} \in \bar{y}+\bar{v}$, while equation 2.6 yields $\bar{v} \in \bar{x}+\bar{z}$. Therefore,

$$
\langle\bar{a}, \overline{v a y}\rangle \equiv\langle\bar{y}, \bar{v}\rangle \text { and }\langle\bar{v}, \overline{v x} \bar{z}\rangle \equiv\langle\bar{x}, \bar{z}\rangle .
$$

These isometries imply that, in order to prove that $\langle\bar{a}, \bar{b}, \bar{c}\rangle \equiv\langle\bar{y}, \bar{x}, \bar{z}\rangle$, it is enough to verify that $\langle\overline{v a y}, \overline{v x z}\rangle \equiv\langle\bar{b}, \bar{c}\rangle$. From the isometries in equation 2.3 we get $\bar{\alpha}=\overline{a x \beta}, \bar{\gamma}=\overline{y z \beta}$ and $\langle\bar{b}, \bar{c}\rangle \equiv\langle\bar{\alpha}, \bar{\gamma}\rangle$. Then, we have $\langle\bar{b}, \bar{c}\rangle \equiv\langle\overline{a x \beta}, \overline{z \beta}\rangle$.
Hence, what is needed is equivalent to $\langle\overline{a x \beta}, \overline{z \beta}\rangle \equiv\langle\overline{v a y}, \overline{v x z}\rangle$. Since the discriminants are
the same, it is enough to prove $\overline{a x \beta} \in \overline{v a y}+\overline{v x z}$.

$$
\overline{a x \beta} \in \overline{v a y}+\overline{v x z} \Leftrightarrow \overline{a x \beta a x v} \in \overline{v a y a x v}+\overline{v x z a x v} \Leftrightarrow \overline{v \beta} \in \overline{x y}+\overline{a z} .
$$

then, it is enough verify that $\overline{v \beta} \in \overline{x y}+\overline{a z}$. Moreover, axiom [DM3], already gave to us that $\bar{v} \bar{\beta} \in \overline{x y}+\overline{a z}$, which finalize the verification of [SG6].

Example 2.3.5. Let $X_{n}$ be the kaleidoscope multiring (as defined in 1.2.12). Of course, if $n \geq 2$, $X_{n}$ is never a DM-hyperfield. However, considering $T=X_{n}^{2} \backslash\{0\}$, since $X_{n}^{2}=\{0,1,2, \ldots, n\}$ we get

$$
K:=X_{n} / m T \cong X_{1}=\{-1,0,1\} .
$$

Since $X_{1}$ is a special hyperfield, $\left(X_{n}, T\right)$ is a DM-multiring.
Example 2.3.6. Let $p$ be a prime integer and consider the $H_{p}$ as defined in 1.2.13 and $T=$ $\sum H_{p}^{2} \backslash\{0\}$. Then $\left(H_{p}, T\right)$ is a DM-hyperfield since $H_{p} / m$ is a real reduced hyperfield.

The above Theorem says that our DM-hyperfields are compatible with the special group structure obtained using Theorem 1.32 of [28].
Theorem 2.3.7. Let $A$ be a domain with $2 \neq 0$. Consider $T \subseteq A$ be a proper preordering or $T=A^{2}$ and denote $T^{*}=T \backslash\{0\}$. Then $A /{ }_{m} T^{*}$ is a special hyperfield, and therefore $G_{T}(A):=\left(A /{ }_{m} T^{*}\right) \backslash\{0\}$ is a special group with representation given by

$$
D_{G_{A}}(\bar{a}, \bar{b})=\bar{a}+\bar{b}=\left\{\bar{c}: c r=a s+b t \text { for some } r, s, t \in T^{*}\right\} .
$$

Moreover, $G_{T}(A)$ is reduced if and only if $T$ is a proper preordering.
Proof. By Theorem 2.3.4. we only need to proof that $A / m T^{*}$ is a DM-hyperfield. First of all, note that

$$
\text { For all } a, b \in A^{*}, \bar{a}, \bar{b} \in \bar{a}+\bar{b} \text {. }
$$

If $a= \pm b$ is immediate (for example, $a(5 a)^{2}=a(4 a)^{2}+a(3 a)^{2}$ or $a(3 a)^{2}=a(5 a)^{2}-a(4 a)^{2}$, in the case where $3,5 \neq 0$ ). If $a \neq \pm b$, then

$$
a(a+b)^{2}=a(a-b)^{2}+b(2 a)^{2}
$$

and $a^{2}+b^{2},(a-b)^{2}, 2 a^{2} \in T^{*}$. Hence $\bar{a} \in \bar{a}+\bar{b}$. Similarly we conclude $\bar{b} \in \bar{a}+\bar{b}$.
Now, we verify the axioms [DM0]-[DM3].
DM0 Of course, $\overline{0} \in \overline{1}-\overline{1}$. If $a \neq 0$, and $a \neq \pm 1$, then

$$
4 a=(a+1)^{2}-(a-1)^{2},
$$

and hence $\bar{a} \in \overline{1}-\overline{1}$. If $a=1$ or $a=-1$, then

$$
9=5^{2}-4^{2} \text { and }-9=4^{2}-5^{2}
$$

testimony that $\overline{1},-\overline{1} \in \overline{1}-\overline{1}$. Therefore $A / m T^{*}$ is hyperbolic.
DM1 Let $\bar{a} \neq 0$ in $A / m$. Then $a^{2} \in T$, hence $\bar{a}^{2}=\overline{1}$.

DM2 Suppose without loss of generality that $a \in A^{*}, a \notin T$ (and hence $\bar{a} \notin\{-\overline{1}, \overline{0}, \overline{1}\}$ ). Now, let $\bar{\alpha}, \bar{\beta} \in \overline{1}+\bar{a}$, with $\alpha x=r+a s, \beta y=t+a w$, for some $x, y, r, s, t, w \in T^{*}$. Then

$$
(r+a s)(t+a w)=\left(r t+a^{2} s w\right)+(s t+r w) a .
$$

If $T$ is a preordering, then $r t+a^{2} s w \in T^{*}$ and $s t+r w \in T^{*}$. If $T=A^{2}$, then $r=r_{1}^{2}, s=s_{1}^{2}$, $t=t_{1}^{2}, w=w_{1}^{2}$ for some $r_{1}, s_{1}, t_{1}, w_{1} \in A^{*}$. Therefore

$$
\begin{aligned}
(r+a s)(t+a w) & =\left(r t+a^{2} s w\right)+(s t+r w) a \\
& =a^{2} s w+r t-2 r_{1} s_{1} t_{1} w_{1} a+2 r_{1} s_{1} t_{1} w_{1} a+(s t+r w) a \\
& =\left(a^{2} s w-2 r_{1} s_{1} t_{1} w_{1} a+r t\right)+\left(s t+2 r_{1} s_{1} t_{1} w_{1}+r w\right) a \\
& =\left(a^{2} s_{1}^{2} w_{1}^{2}-2 r_{1} s_{1} t_{1} w_{1} a+r_{1}^{2} t_{1}^{2}\right)+\left(s_{1}^{2} t_{1}^{2}+2 r_{1} s_{1} t_{1} w_{1}+r_{1}^{2} w_{1}^{2}\right) a \\
& =\left(a s_{1} w_{1}-r_{1} t_{1}\right)^{2}+\left(s_{1} t_{1}+r_{1} w_{1}\right)^{2} a .
\end{aligned}
$$

If $\left(a s_{1} w_{1}-r_{1} t_{1}\right)^{2}=\left(s_{1} t_{1}+r_{1} w_{1}\right)^{2}=0$ we have $\overline{r+a t}=\overline{0}$ or $\overline{s+a w}=\overline{0}$, and hence $r=-a t$ or $s=-a w$, and both cases imply $-\bar{a}=\overline{1}$. If $\left(a s_{1} w_{1}-r_{1} t_{1}\right)^{2},\left(s_{1} t_{1}+r_{1} w_{1}\right)^{2} \neq 0$ then $\left(a s_{1} w_{1}-r_{1} t_{1}\right)^{2},\left(s_{1} t_{1}+r_{1} w_{1}\right)^{2} \in T^{*}$ and we are done. If $\left(a s_{1} w_{1}-r_{1} t_{1}\right)^{2}=0$, using 2.8

$$
(r+a s)(t+a w)=\left(s_{1} t_{1}+r_{1} w_{1}\right)^{2} a \Rightarrow \bar{\alpha} \bar{\beta}=\bar{a} \in 1+\bar{a}
$$

If $\left(s_{1} t_{1}+r_{1} w_{1}\right)^{2}=0$, using 2.8

$$
(r+a s)(t+a w)=\left(a s_{1} w_{1}-r_{1} t_{1}\right)^{2} \Rightarrow \bar{\alpha} \bar{\beta}=\overline{1} \in 1+\bar{a}
$$

completing the proof.

DM3 Let

$$
\left\{\begin{array}{l}
\bar{a} \in \bar{x}+\bar{b} \\
\bar{b} \in \bar{y}+\bar{z}
\end{array} \quad \text { in } A / m T,\right.
$$

with $\bar{a}, \bar{b}, \bar{x}, \bar{y}, \bar{z} \neq \overline{0}$. Then, there exists $p_{a}, q_{a}, r_{a}, p_{b}, q_{b}, r_{b} \in T$ such that

$$
\begin{gather*}
a p_{a}=x q_{a}+b r_{a} .  \tag{eqz}\\
b p_{b}=y q_{b}+z r_{b} . \tag{}
\end{gather*}
$$

Therefore

$$
a p_{a} p_{b}=x p_{b} q_{a}+b p_{b} r_{a}=x p_{b} q_{a}+\left(y q_{b}+z r_{b}\right) r_{a}=x p_{b} q_{a}+y q_{b} r_{a}+z r_{a} r_{b} .
$$

Now, consider

$$
v=x p_{b} q_{a}+z r_{a} r_{b} .
$$

Note that $\bar{v} \in \bar{x}+\bar{z}$ and

$$
\begin{equation*}
a p_{a} p_{b}=y q_{b} r_{a}+v, \tag{5}
\end{equation*}
$$

with $\bar{a} \in \bar{y}+\bar{v}$. In order to complete the proof, we only need to verify that $\overline{v b} \in \overline{x y}+\overline{a z}$. In
fact,

$$
\begin{aligned}
v b p_{b} & =\left(x p_{b} q_{a}+z r_{a} r_{b}\right)\left(y q_{b}+z r_{b}\right) \\
& =x y p_{b} q_{a} q_{b}+x z p_{b} q_{a} r_{b}+y z q_{b} r_{a} r_{b}+z^{2} r_{a} r_{b}^{2} \\
& =x y p_{b} q_{a} q_{b}+z\left(x p_{b} q_{a} r_{b}+y q_{b} r_{a} r_{b}+z r_{a} r_{b}^{2}\right) \\
& =x y p_{b} q_{a} q_{b}+z\left(x p_{b} q_{a} r_{b}+\left(y q_{b}+z r_{b}\right) r_{a} r_{b}\right) \\
& =x y p_{b} q_{a} q_{b}+z\left(x p_{b} q_{a} r_{b}+b p_{b} r_{a} r_{b}\right) \\
& =x y p_{b} q_{a} q_{b}+\left(x q_{a}+b r_{a}\right) z p_{b} r_{b} \\
& =x y p_{b} q_{a} q_{b}+a p_{a} z p_{b} r_{b} \\
& =x y p_{b} q_{a} q_{b}+a z p_{a} p_{b} r_{b}
\end{aligned}
$$

and hence, $\overline{v b} \in \overline{x y}+\overline{a z}$.

Corollary 2.3.8. Let $D$ be a domain with $2 \neq 0$ and consider the polynomial ring $D\left[x_{1}, \ldots, x_{n}\right]$. Let $T \subseteq D\left[x_{1}, \ldots, x_{n}\right]$ be a preordering or $T=\left(D\left[x_{1}, \ldots, x_{n}\right]\right)^{2}$. Then $D\left[x_{1}, \ldots, x_{n}\right] / m T^{*}$ is a special group.

Theorem 2.3.9. Let $F$ be a hyperfield satisfying DM0-DM2. Then $F$ satisfy DM3 if and only if satisfy SMF4. In other words, $F$ is a DM-hyperfield if and only if is a special hyperfield.

Proof. After Theorem 2.3.4, we only need to prove that if $F$ is a special hyperfield then $F$ satisfy DM3. Let $a \in x+b$ and $b \in y+z$. Then by definition, $a \in x+y+z$, and then, there exist some $v \in x+z$ such that $a \in y+v$. We need to prove that $v b \in x y+a z$. We discuss two cases.

Case I: $v=0$. Then $a=y$ and $z=-x$. Moreover

$$
0=v b \in a x-a x=x y+a z
$$

Case II: $v \neq 0$. Here we consider the special group structure in $F^{*}$. Moreover, for all $a, b \in F^{*}$, $a, b \in a+b$. Considering $a \in x+b$ and $b \in y+z$, we get the above isometries

$$
\langle b y z, x\rangle \equiv\langle x, b y z\rangle,\langle a x b, a\rangle \equiv\langle x, b\rangle \text { and }\langle y, z\rangle \equiv\langle b y z, b\rangle
$$

Then by definition $\langle b y z, a x b, a\rangle \equiv\langle x, y, z\rangle$.
Moreover, considering $a \in y+v$ and $v \in x+z$, we get the above isometries

$$
\langle v x z, y\rangle \equiv\langle y, v x z\rangle,\langle a y v, a\rangle \equiv\langle y, v\rangle \text { and }\langle x, z\rangle \equiv\langle v x z, v\rangle
$$

Then by definition $\langle v x z, a y v, a\rangle \equiv\langle y, x, z\rangle$. Since $F^{*}$ is a special group, $\langle x, y, z\rangle \equiv\langle y, x, z\rangle$ and the isometry relation is 3 -transitive. Then

$$
\langle b y z, a x b, a\rangle \equiv\langle x, y, z\rangle \equiv\langle y, x, z\rangle \equiv\langle v x z, a y v, a\rangle
$$

and hence, $\langle b y z, a x b, a\rangle \equiv\langle v x z, a y v, a\rangle$. Using Witt's Cancellation, $\langle b y z, a x b\rangle \equiv\langle v x z, a y v\rangle$. Then,

$$
v x z \in b y z+a x b \Rightarrow v b x z \in y z+a x \Rightarrow v b \in x y+a z
$$

completing the proof.

Theorem 2.3.10. Let $(G, \equiv, 1,-1)$ be a pre-special group. Are equivalent:

1. $G$ is special, i.e, satisfy (for example) SG6.
2. $M(G)$ (the hyperfield associated to $G$ ) satisfy DM3.
3. $G$ satisfy the following condition for all $a, b, x, y, z \in G$ :

$$
\begin{aligned}
& \text { If } a \in D_{G}(x, b) \text { and } b \in D_{G}(y, z) \text { then there exist } v \in D_{G}(x, z) \\
& \text { such that } a \in D_{G}(y, v) \text { and } v b \in D_{G}(x y, a z) .
\end{aligned}
$$

### 2.4 DM-multirings and Quadratically presentable fields

Let $(R, T)$ be a DM-multiring and $G(R, T):=(R / m T) \backslash\{0\}$. Since $G(R, T)$ is a special group, we can provide a theory of quadratic forms for $R$ inherited from $G(R, T)$ : Let $\equiv$ be the isometry relation on $G(R, T)^{2}$ given by $\langle a, b\rangle \equiv\langle c, d\rangle$ iff $a b=c d$ in $G(R, T)$ and $a \in c+d \backslash\{0\}$. We extend $\equiv$ to a binary relation $\equiv_{n}$ on $G(R, T)^{n}$, by induction on $n \geq 2$, as follows:

$$
\text { i }-\equiv_{2}=\equiv .
$$

ii - $\left\langle a_{1}, \ldots, a_{n}\right\rangle \equiv_{n}\left\langle b_{1}, \ldots, b_{n}\right\rangle$ if and only there are $x, y, z_{3}, \ldots, z_{n} \in A$ such that $\left\langle a_{1}, x\right\rangle \equiv\left\langle b_{1}, y\right\rangle$, $\left\langle a_{2}, \ldots, a_{n}\right\rangle \equiv_{n-1}\left\langle x, z_{3}, \ldots, z_{n}\right\rangle$ and $\left\langle b_{2}, \ldots, b_{n}\right\rangle \equiv_{n-1}\left\langle y, z_{3}, \ldots, z_{n}\right\rangle$.

Since $G(R, T)$ is a special group, $\equiv_{n}$ is transitive for all $n \geq 2$ (in fact, this is the content of axiom SG6). Whenever clear from the context, we frequently abuse notation and indicate the aforedescribed extension $\equiv$ by the same symbol.

A form $\varphi$ on $G(R, T)$ is an $n$-tuple $\left\langle a_{1}, \ldots, a_{n}\right\rangle$ of elements of $G(R, T) ; n$ is called the dimension of $\varphi, \operatorname{dim}(\varphi)$. We also call $\varphi$ a $n$-form.

By convention, two forms of dimension 1 are isometric if and only if they have the same coefficients. If $\varphi=\left\langle a_{1}, \ldots, a_{n}\right\rangle$ is a form on $G(R, T)$, define
a - The set of elements represented by $\varphi$ as

$$
D_{G(R, T)}(\varphi):=\left\{b \in G(R, T): \exists z_{2}, \ldots, z_{n} \in G(R, T) \text { such that } \varphi \equiv\left\langle b, z_{2}, \ldots, z_{n}\right\rangle\right\}
$$

b - The discriminant of $\varphi$ as $d(\varphi)=\prod_{i=1}^{n} a_{i}$.
c- Direct sum as $\varphi \oplus \theta=\left\langle a_{1}, \ldots, a_{n}, b_{1}, \ldots, b_{m}\right\rangle$.
d - Tensor product as $\varphi \otimes \theta=\left\langle a_{1} b_{1}, \ldots, a_{i} b_{j}, \ldots, a_{n} b_{m}\right\rangle$. If $a \in G(R, T),\langle a\rangle \otimes \varphi$ is written $a \varphi$.
A form $\varphi$ on $G(R, T)$ is isotropic if there is a form $\psi$ over $G(R, T)$ such that $\varphi \equiv\langle 1,-1\rangle \oplus \psi$; otherwise it is said to be anisotropic. We say that $\varphi$ is universal if $D_{G(R, T)}(\varphi)=G(R, T)$.

In this sense, Witt Ring $W(R, T)$ of $(R, T)$ is defined as the Witt ring $W(G(R, T))$ of $G(R, T)$. We can go further, and define a form $\varphi=\left\langle a_{1}, \ldots, a_{n}\right\rangle$ on $(R, T)$ by considering the form $\bar{\varphi}:=$ $\left\langle\overline{a_{1}}, \ldots, \overline{a_{n}}\right\rangle$ on $G(R, T)$ and so on.

Moreover, this quadratic form theory inherited from $G(R, T)$ is compatible with the more general Witt rings described by P. Gladik and K. Worytkiewicz in 37]:

Definition 2.4.1 (Presentable monoid, group, ring [37]). Let $(A, \leq, 0)$ be a pointed poset (i.e, a poset with a distinguished element $0 \in A$ ).
$a-(A, \leq, 0,+)$ is a presentable monoid if the distinguished element 0 is supercompact and + : $M \times M \rightarrow M$ is a suprema-preserving binary operation such that for all $a, b, c \in M$
(a) $a+(b+c)=(a+b)+c$;
(b) $a+0=0+a=a$;
(c) $a+b=b+a$.
$b-(A, \leq, 0,+,-)$ is a presentable group if $(A, \leq, 0,+)$ is a presentable monoid and $-: G \rightarrow G$ is a suprema preserving involutive homomorphism (called inversion) such that $s \leq t+u$ imply $t \leq s+(-u)$ for all $s, t, u \in \mathcal{S}_{G}$ (here $\mathcal{S}_{G}$ denote the set of $G$ 's minimal elements).
$c-(A, \leq, 0,1,+,-, \cdot)$ is a presentable ring if $(A, \leq, 0,+,-)$ is a presentable group, $(A, 1, \cdot)$ is a commutative monoid such that the element 1 is supercompact, • is compatible with $\leq$ and (i.e, $a \leq b$ imply $a \cdot c \leq b \cdot c$ and $a \cdot(-b)=-(a \cdot b)$ for all $a, b, c \in A$ ), is distributive with respect to $+0 \cdot a=0$ for all $a \in R$ and $\cdot$ satisfy

$$
\mathcal{S}_{a \cdot b}=\left\{s \cdot t: s \in \mathcal{S}_{a}, t \in \mathcal{S}_{b}\right\} .
$$

Here $\mathcal{S}_{a}:=\downarrow a \cap \mathcal{S}_{A}$ for all $a \in A$, i.e, $\mathcal{S}_{a}$ is the set of all minimal elements below $a \in A$.
$d-(A, \leq, 0,1,+,-, \cdot)$ is a presentable field if is a presentable ring such that every non-zero element is invertible.

Now we recall the concept of quadratically presentable fields (in the sense of Definitions 5.1, 5.5 and 5.7 of [37]). A presentable field ( $A, \leq, 0,1,+,-, \cdot)$ is pre-quadratically presentable whenever
i- $a \leq a+b$ for all $a \in \mathcal{S}_{A}^{*}, b \in \mathcal{S}_{A}$;
ii - $a \leq 1+b$ and $a \leq 1+c$ imply $a \leq 1-b c$ for all $a, b, c \in \mathcal{S}_{A}$;
iii - $a^{2}=1$ for all $a \in \mathcal{S}_{A} \backslash\{0\}$.
A form on a pre-quadratically presentable field $A$ is an $n$-tuple $\left\langle a_{1}, \ldots, a_{n}\right\rangle$ of elements of $\mathcal{S}_{A}^{*}$. The relation $\cong$ of isometry of forms of the same dimension is given by induction: (i) $\langle a\rangle \cong\langle b\rangle$ iff $a=b$; (ii) $\left\langle a_{1}, a_{2}\right\rangle \cong\left\langle b_{1}, b_{2}\right\rangle$ iff $a_{1} a_{2}=b_{1} b_{2}$ and $b_{1} \leq a_{1}+a_{2}$; (iii) finally, for $n \geq 3$

$$
\begin{gathered}
\left\langle a_{1}, \ldots, a_{n}\right\rangle \cong\left\langle b_{1}, \ldots, b_{n}\right\rangle \text { iff there exists } x, y, c_{3}, \ldots, c_{n} \in \mathcal{S}_{A}^{*} \text { such that }\left\langle a_{1}, x\right\rangle \cong\left\langle b_{1}, y\right\rangle \\
\left\langle a_{2}, \ldots, a_{n}\right\rangle \cong\left\langle x, c_{3}, \ldots, c_{n}\right\rangle,\left\langle a_{b}, \ldots, b_{n}\right\rangle \cong\left\langle y, c_{3}, \ldots, c_{n}\right\rangle .
\end{gathered}
$$

A pre-quadratically presentable field is quadratically presentable whenever the isometry relation defined above is an equivalence relation on the set of all forms of the same dimension.

Let $(R, T)$ be a DM-multiring. Let $K:=R /{ }_{m} T$ and consider $\mathcal{P}^{*}(K)$, the pierced powerset of the set $K$ (that is, its set of nonempty subsets). Then $\left(\mathcal{P}^{*}(K), \subseteq,\{0\},\{1\},+,-, \cdot\right)$ is a presentable field ([37], Example 4.5), where the operations in $\mathcal{P}^{*}(K)$ are defined for $A, B \in \mathcal{P}^{*}(K)$ by

$$
-A:=\bigcup_{a \in A}\{-a\}, A+B:=\bigcup_{a \in A, b \in B} a+b \text { and } A \cdot B:=\bigcup_{a \in A, b \in B}\{a \cdot b\} .
$$

Following 5.18 [37], we obtain:

Theorem 2.4.2. Let $(R, T)$ be a $D M$-multiring. Let $K:=R / m T$ and $\left(\mathcal{P}^{*}(K), \subseteq,\{0\},\{1\},+,-, \cdot\right)$ be the induced presentable field. Then:

1. $\mathcal{P}^{*}(K)$ is a quadratically presentable field.
2. $W\left(\mathcal{P}^{*}(K)\right) \cong W(K)=W(R, T)$, where $W\left(\mathcal{P}^{*}(K)\right)$ is the Witt ring defined in 5.13 37].

Proof. (1) This follows, essentially, from the same argument of 2.3 .4 , since $K$ is a special hyperfield.
(2) Just repeat the arguments used in 7.1, 7.2 and 7.3 of [37].

For the readers comfortable with theory of special groups, the proof of this Theorem is just a translation of axiom SG6.

In 7.4 of 37 is asked:
"It is an open question when the resulting pre-quadratically presentable field is quadratically presentable."

We finish this section arguing that such question is, in principle, non void. More precisely:

## qgladik

Proposition 2.4.3. There exists a pre-quadratically presentable field that is not quadratically presentable.

Proof. We show that $p Q P F$ is a cocomplete category but $Q P F$ is not a cocomplete category.

- In 5.18 of [37] are established equivalences of categories:
quadratically presentable fields $(Q P F) \nrightarrow s$ special groups $(S G)$;
pre-quadratically presentable fields $(p Q P F) \nrightarrow$ pre-special groups $(p S G)$.
- $p Q P F(\simeq p S G)$ is a cocomplete category. According the Definition of pre-special group (Definition 1.2 in [28]), it is axiomatized by a universal Horn Theory (Definition 5.10 in [1]) thus it is a limit theory (Definition 5.7 in [1]). By Theorem 5.9 in [1], $p S G$ is a finitely locally presentable category, (Definition 1.9 in [1), thus it is a cocomplete category.
- $Q P F(\simeq S G)$ is not a cocomplete category.
* Consider $R S G$ the full subcategory of $S G$ of all reduced special groups, i.e. a special group $G$ such that for each $a \in G,\langle a, a\rangle \equiv\langle 1,1\rangle$ iff $a=1$. This is a slightly variation on the notion of reduced special group (Definition 1.2 in [28]) since we not exclude the case where $G=\{1\}$ (equivalently, we not impose $-1 \neq 1$ ). Following the proofs of the results in Chapter 10 , Section 3, in [28], the category $R S G$ of all reduced special groups (including the trivial special group $\{1\})$ ) misses some binary coproducts, thus is not cocomplete (see for instance, Proposition 10.11 of [28]).
* The full subcategory $\iota: R S G \hookrightarrow S G$ is reflexive, i.e. it has a left adjoint $S: S G \rightarrow R S G$, $G \in \operatorname{Obj}(S G) \mapsto G / \operatorname{Sat}(G) \in \operatorname{Obj}(R S G)$, where the unity of adjunction is $\left(G \xrightarrow{q_{G}} S(G):=\right.$ $G / \operatorname{Sat}(G))_{G \in O b j(S G)}$. This follows from a combination of results in [28]: Remark (iii) just below Definition 2.7; Remark 2.16 and Proposition 2.21.
* Let $\Gamma: \mathcal{I} \rightarrow R S G$ be a small diagram that does not have a colimit in $R S G$. Suppose that $\iota \circ \Gamma: \mathcal{I} \rightarrow S G$ has a colimit $\left(\gamma_{i}: \Gamma(i) \rightarrow G_{\infty}\right)_{i \in o b j(\mathcal{I})}$ in $S G$. Then it is easy to check that $\left(q_{G_{\infty}} \circ \gamma_{i}: \Gamma(i) \rightarrow S\left(G_{\infty}\right)\right)_{i \in O b j(\mathcal{I})}$ satisfies the universal property of being the colimit of $\Gamma: \mathcal{I} \rightarrow R S G$ in $R S G$, a contradiction.


### 2.5 Quadratic Multirings and (Formally) Real Semigroup associated to Semi real rings via Marshall quotient

Paraphrasing M. Marshall, "when we change fields for rings, we are in deep water" (48))! For example, let $R$ be a generic commutative ring and $T \subseteq R$ be a multiplicative set containing 1 . From now on, we denote

$$
\begin{aligned}
z d(R) & :=\{a \in R: a \text { is a zero divisor }\} \\
n z d(R) & :=R \backslash z d(R)=\{a \in R: a \text { is not a zero divisor }\} .
\end{aligned}
$$

If $a, b \in T \backslash\{0\}$ with $a b=0$ (i.e, $a, b$ are zero-divisors), then $R / m T^{*} \cong\{0\}$ : in fact for all $x \in R, x(a b)=0 \cdot 1$ with $a b, 1 \in T$, and hence $\bar{x}=\overline{0}$. Even in the case $T \subseteq n z d(R)$, if $a \in z d(R)$, say $a b=0$ for some $b \in z d(R)$ then $\bar{a} \bar{b}=\overline{0}$, so $(\overline{a b})^{2}=0 \neq \overline{1}$, and in particular, $R / m T$ is not a hyperfield.

Then, if $z d(R) \neq \emptyset, R / m T^{*}$ will never be a special group, since will never be a hyperfield. Because this, we seek for conditions for a pair $(R, T)$ with $R$ a ring and $T \subseteq n z d(R)$ multiplicative provide a (formally) real semigroup structure in $R / m T$.

In this context we christen the following Definition:

## qring

Definition 2.5.1. Let $R$ be a multiring and $T \subseteq n z d(R)$ be a multiplicative subset containing 1 . We say that $(R, T)$ is a quadratic pair if
Q1 $R /{ }_{m} T$ is semi real.
Q2 If $a \in R$ and $a^{2} \notin z d(R)$, then $a^{2} \in T$.
Q3 For all $a \in R$, then $\bar{a}^{3}=\bar{a}$ in $R /{ }_{m} T$.
Q4 For all $a, b \in R$, there exists $r, s, t \in T$ such that $a r \in a^{3} s+a^{2} b t$.
Let's look closely to the axioms in Definition 2.5.1. In this sense, Q1 is a kind of generalization of the semireal condition and Q 2 is a weakness of $A^{2} \subseteq T$. The following Theorem is immediate:
qringteo
Theorem 2.5.2. Let $(R, T)$ be a quadratic pair and define for all $a, b, c \in R$ the following relations:

$$
\begin{array}{r}
\bar{c} \in D^{t}(\bar{a}, \bar{b}) \text { if and only if } \bar{c} \in \bar{a}+\bar{b} \\
\bar{c} \in D(\bar{a}, \bar{b}) \text { if and only if } \bar{c} \in D^{t}\left(\bar{c}^{2} a, \bar{c}^{2} b\right) .
\end{array}
$$

Then $\left(R /{ }_{m} T, D, D^{t}\right)$ is a formally real semigroup. Conversely, if $\left(G, D, D^{t}\right)$ is a formally real semigroup such that $a^{2}$ is a zero divisor or $a^{2}=1$. Define

$$
c \in a+b \text { if and only if } c \in D^{t}(a, b) .
$$

Then $(G,\{1\})$ is a quadratic pair.
Proof. Let $(R, T)$ be a quadratic pair. Axiom RS7b is consequence of Q1 and axiom RS1 is consequence of Q4. The other axioms of formally realsemigroup are consequence of basic properties of multiring and so on.

Conversely, if $\left(G, D, D^{t}\right)$ is a formally real semigroup such that $a^{2}$ is a zero divisor or $a^{2}=1$, we automatically have Q2. Q1 is consequence of RS7b, Q3 is consequence of $G$ be a ternary semigroup and Q4 is consequence of RS1. The fact of $(G,+, \cdot, 0,1)$ be a multiring is consequence of the another axioms of formally realsemigroup (and ternary semigroup).

### 2.5. QUADRATIC MULTIRINGS AND (FORMALLY) REAL SEMIGROUP ASSOCIATED TO SEMI REAL RI

Now is time to deal with the real semigroup case. We define the following:
Definition 2.5.3. A Dickmann-Petrovich multiring (or DP-multiring for short) ${ }^{2}$ is a quadratic pair $(R, T)$ satisfy the following properties:

DP1 $1+T \subseteq T$.
DP2 For all $a \in R$, exist $t \in T$ such that $1+a^{2} t \in T$.
DP3 For all $a, b \in R, \bar{a}^{2}+\bar{b}^{2}$ is a singleton set in $R /{ }_{m} T$.
teopmr
Theorem 2.5.4. Let $(R, T)$ be a DP-ring and denote $R s(R)=(R / m T)$. Then $R s(R)$ is a real reduced multiring (thus it is a real semigroup).

Proof. Since $T \subseteq n z d(R), \overline{1} \neq \overline{0}$ in $\operatorname{Rs}(R)$. Moreover, by (Q4) we get $\overline{a^{3}}=\bar{a}$ in $R s(R)$.
Note that since $T$ is multiplicative, [Q0] and [DP1] imply $T \cdot T=T$ and

$$
T+T=T+T \cdot T=T \cdot(1+T) \subseteq T \cdot T=T
$$

then we have that $T+T \subseteq T$ which imply that $\bar{a}+\bar{a}=\{\bar{a}\}$ for all $\bar{a} \in R s(R)$.
From (DP2) we get $\overline{1}+\bar{b}^{2}=\{\overline{1}\}$ for all $b \in R$, which imply $\bar{a}+\overline{a b^{2}}=\{\bar{a}\}$ for all $a, b \in R$. Finally, [DP3] says that $\bar{a}^{2}+\bar{b}^{2}$ is a singleton set in $R / m T$, completing the proof that $R / m T$ is a real semigroup.

Example 2.5.5. Let $(R, T)$ be a DM-multiring. Then $(R, T)$ is also a quadratic pair.
Example 2.5.6. Let $(R, T)$ be a $D M$-ring such that $T+T \subseteq T$. Then $(R, T)$ is also a $D P$-ring.
With Definition 2.5.1 and Theorem 2.5.2, we generalize the real reduced multirings:
quadraticring
Definition 2.5.7. A multiring $A$ is said to be quadratic if satisfy the following properties:
QM0 $-1 \notin \sum A^{2}$.
QM1 for all $a \in A, a \in 1-1$.
QM2 for all $a \in A, a^{3}=a$.
QM3 for all $a, b \in A, a \in a+a^{2} b$.
Example 2.5.8. Let $p$ be a prime integer and consider $H_{p}$ as in 1.2.13. Since $a^{2}=1$ and $a=-a$ for all $a \in H_{p}$ and $a+a=H_{p}$ for all $a \neq 0$, we have that $H_{p}$ is not a quadratic multiring.

But $H_{p}$ satisfy QM1, QM2 and QM3. Then, consider the product multiring $R=X_{1} \times H_{p}$, where $X_{1}=\{-1,0,1\}$. Since $X_{1}$ is a DM-hyperfield (and hence a DP-multiring) and the operations and multioperation in $R$ is defined coordinatewise, we have that $R$ satisfy QM1, QM2 and QM3. Since $(a, b) \in R^{2}$ if and only if $a \in\{0,1\}$ and $b \in H_{p}$, we have $-1_{R}=(-1,1) \notin R^{2}$. Hence $R$ is $a$ quadratic multiring.

Example 2.5.9 (Constructions).

[^5]$i$ - (Products) Let $\left\{R_{i}\right\}_{i \in I}$ be a class of quadratic multiring and let $R=\prod_{i \in I} R_{i}$. Since the operations and multioperation in $R$ is defined coordinatewise, we have that $R$ is a quadratic multiring. More generally, suppose that $R_{i}$ satisfy $Q M 1, Q M 2$ and $Q M 3$ for all $i \in I$. If there is an index $i_{0} \in I$ such that $R_{i_{0}}$ is a quadratic multiring, then $R$ is a quadratic multiring.
ii - (Directed Colimits) If $(I, \leq)$ is an upward directed poset and $\left(f_{i j}: R_{i} \rightarrow R_{j}\right)_{i \leq j}$ is a diagram of quadratic multirings, then colim $i_{i \in I} R_{i}$ is a quadratic multiring. More generally, if ( $f_{i j}$ : $\left.R_{i} \rightarrow R_{j}\right)_{i \leq j}$ is an upward directed diagram of multirings such that $\left\{i \in I: R_{i}\right.$ is a quadratic multiring\} is a cofinal subset of $I$, then colim ${ }_{i \in I} R_{i}$ is a quadratic multiring.
iii - (Reduced Products and Ultraproducts) The class of quadratic multirings can be axiomatized by certain kind of first-order formulas (in a convenient relational language) that shows that this subclass of the class of multirings is closed under reduced products (and ultraproducts, in particular). This result can be achieved more directly by the description of reduced product of a family of (quadratic) multirings, modulo some filter on the index set, as the directed colimit of products of the members of the family indexed by some member of the filter: $\prod_{i \in I} R_{i} / \mathcal{F} \cong$ $\operatorname{colim}_{J \in \mathcal{F}} \prod_{i \in J} R_{i}$.
sgqring
Example 2.5.10 (Special Groups). Let $G$ be a special group, and consider $F=M(G):=G \cup\{0\}$ its special hyperfield associated. Of course, $F$ satisfy conditions QM1-QM3 in 2.5.7. F satisfy DM0 iff $F$ is formally real, i.e, if $-1 \notin \sum F^{2}$, which occurs iff $G$ is formally real, i.e,
$$
-1 \notin D_{G}(n \otimes\langle 1\rangle) \text { for all } n \geq 1 .
$$

Example 2.5.11. Let $A$ be a von Neumann regular semi-real ring such that $2 \in A^{\times}$. Then $A /{ }_{m} A^{\times 2}$ is a quadratic multiring. In fact, first observe that
i) If $F$ is a field with $2 \in F^{\times}$, then $F /{ }_{m} F^{\times 2}$ is a multiring that satisfies $\mathbf{Q M 1 - Q M 3}$ as indicate Examples 2.1.4 and 2.5.10. This means that $F$ satisfies the following Horn-geometric sentences:

- $\forall a \exists x, y, x^{\prime}, y^{\prime}\left(x x^{\prime}=y y^{\prime}=1 \wedge a=x^{2}-y^{2}\right)$.
- $\forall a \exists x, y, x^{\prime}, y^{\prime}\left(x x^{\prime}=y y^{\prime}=1 \wedge a^{3} x^{2}=a y^{2}\right)$.
- $\forall a, b \exists x, y, z, x^{\prime}, y^{\prime}, z^{\prime}\left(x x^{\prime}=y y^{\prime}=z z^{\prime}=1 \wedge a x^{2}=a y^{2}+a^{2} b z^{2}\right)$.
ii) The Proposition 5.6 of [31] shows that the von Neumann regular ring $A$ is the ring of global sections over a Boolean space where the sheaf has fields with 2 invertible as stalks.
Thus, the Proposition 3.2-(d), 31], applied to the sheaf of item ii) above implies that formulas of item $i$ ) are valid in $A$. Therefore $A / m A^{\times 2}$ is a quadratic multiring.
Example 2.5.12 (Faithfully Quadratic Rings). Now, we relate our DM-multirings, DP-multirings and quadratic multirings with faithfully quadratic rings as presented in [32]: let $A$ be a semi-real ring with $2 \in \dot{A}, T$ be a preordering of $A$ or $T=A^{2}$. A T-subgroup of $A$ is a multiplicative subset $S$ of $\dot{A}$ containing $\{-1\} \cup \dot{T}$. For $a, b \in S$, denote

$$
\begin{aligned}
& D_{S, T}^{v}(a, b):=\{c \in S: c=a s+b t \text { for some } s, t \in T\} . \\
& D_{S, T}^{t}(a, b):=\{c \in S: c=a s+b t \text { for some } s, t \in \dot{T}\} .
\end{aligned}
$$

The triple $(A, T, S)$ is faithfully quadratic if (among other things) satisfy $D_{S, T}^{v}(a, b)=D_{S, T}^{t}(a, b)$ for all $a, b \in S$ (see for instance, Definition 3.1 in [32]). Denote

$$
a^{T}=b^{T} \text { iff } a b \in \dot{T} \text { iff } b=a t \text { for some } t,
$$

### 2.5. QUADRATIC MULTIRINGS AND (FORMALLY) REAL SEMIGROUP ASSOCIATED TO SEMI REAL RI

and consider $G_{T}(S)=\left\{a^{T}: a \in S\right\}$. Define the binary isometry $\equiv_{T}^{S}$ by

$$
\left\langle a^{T}, b^{T}\right\rangle \equiv\left\langle c^{T}, d^{T}\right\rangle \text { iff } a^{T} b^{T}=c^{T} d^{T} \text { and } D_{S, T}^{v}(a, b)=D_{S, T}^{v}(c, d) \text {. }
$$

In general, $\left(G_{T}(S), \equiv_{T}^{S},-1^{T}\right)$ is a proto-special group. If $(A, T, S)$ is faithfully quadratic, then Dickmann and Miraglia showed (see Theorem 3.5[32]) that $G_{T}(S)$ is a special group.

Now, consider $(A, T, S)$ and let $R=A / m(T \cap n z d(A))$. Then $D_{S, T}^{t}(a, b) \subseteq \bar{a}+\bar{b}$ for all $a, b \in A$. Moreover, if $A^{2} \subseteq n z d(A)$, or more generally, if $(A, T)$ is a quadratic ring, then $R$ is a quadratic multiring containing the proto special group $G_{T}(S)$. This is particularly useful given that $(A, T, S)$ is not necessarily faithfully quadratic.

Definition 2.5.13. Let $(X, \tau)$ be a topological space. The topology $\tau$ is called perfectly normal if it is normal and every closed set is $G_{\delta}$-set. The topology $\tau$ is called $T_{6}$ if it is Hausdorff and perfectly normal.

## Example 2.5.14.

- A $T_{1}$ topological space $X$ is perfectly normal if, and only if, for every closed set $F$ exists a continuous function $f: X \rightarrow \mathbb{R}$ such that $F=f^{-1}(0)$ (Theorem 1.5.19 of [35]).
- Every metric space is $T_{6}$ (Corollary 4.1.13 of [35]).

Example 2.5.15 (The ring of continuous functions). Let $X$ be $T_{6}$ topological space and consider $A=C(X, \mathbb{R})$, the ring of continuous functions $f: X \rightarrow \mathbb{R}$. Let $T=A^{2} \cap n z d(A)$. In the following, is proved that $C(X, \mathbb{R}) / m T$ is a real reduced multiring (in particular, a quadratic multiring). Before that, consider the remarks:

- Since $X$ is perfectly normal, given a open set $U \subseteq X$ there is a continuous function $g: X \rightarrow \mathbb{R}$ such that $\left.g\right|_{U}$ is strictly positive and $Z(g)=U^{c}$.
- $f \in C(X, \mathbb{R})$ is zero divisor if, and only if, $Z(f)$ has non-empty interior. In fact, if $U \subseteq Z(f)$ is non-empty interior, then exists $g \in C(X, \mathbb{R})$ such that $Z(g)=U^{c}$; thus $g$ is a non-zero function and $f g=0$. Reciprocally, if $Z(f)$ has empty interior and $g \in C(X, \mathbb{R})$ satisfies $f g=0$, then $Z(f)^{c}$ is open and dense while $Z(f)^{c} \subseteq Z(g)$. Since $g$ is continuous, $g=0$ and so $f$ is non-zero divisor.
- By the preceding item,

$$
T=\{f \in C(X, \mathbb{R}): f \text { is non-negative and } Z(f) \text { has empty interior }\} .
$$

Before proceeding with the proof, a notation: given $h \in C(X, \mathbb{R})$, denote by $p_{h} \in C(X, \mathbb{R})$ any function satisfying:

- $Z\left(p_{h}\right)$ has empty interior (i.e. $p_{h}$ is a non-zero divisor).
- $p_{h}$ is non-negative over $Z(h)$.
- For all $x \notin Z(h), p_{h}(x)=h(x)$.

A possible construction is to consider a positive function $p \in C(X, \mathbb{R})$ with $Z(p)=(\operatorname{int}(Z(h)))^{c}$ and set $p_{h}:=h+p$.

Claim. Let $f, g \in C(X, \mathbb{R})$ be two functions and $D \subseteq X$ a dense subset such that for all $x \in D$, $\operatorname{sgn}(f(x))=\operatorname{sgn}(g(x))$. Then $\bar{f}=\bar{g}$ in $C(X, \mathbb{R}) / m T$.

Proof. Assume that for all $x \in D, \operatorname{sgn}(f(x))=\operatorname{sgn}(g(x))$. Then for all $x \in D$ we have $f(x) \cdot p_{|g|}(x)=$ $g(x) \cdot p_{|f|}(x)(*)$. Since $D$ is a dense subset of $X$, the equality $(*)$ is true for all real number. Thus, since $p_{|f|}, p_{|g|} \in T$, we have $\bar{f}=\bar{g}$ in $A /{ }_{m} T$.

To finalize this Example, we have to prove the axioms of real reduced multiring:

- Since $0 \notin T$, we have $\overline{1} \neq \overline{0}$ in $A / m T$.
- For all $x \in \mathbb{R}$, we have $\operatorname{sgn}\left(f^{3}(x)\right)=\operatorname{sgn}(f(x))$. Thus by the above claim $\bar{f}^{3}=\bar{f}$ in $A / m T$.
- Let $f, g \in A$ and $\bar{h} \in \bar{f}+\bar{f} \bar{g}^{2}$ in $A / m$. Then exists $s_{1}, s_{2}, s_{3} \in T$ such that $h s_{1}=f s_{2}+f g^{2} s_{3}$. Thus, for all $x \in Z\left(s_{1}\right)^{c} \cap Z\left(s_{2}\right)^{c} \cap Z\left(s_{3}\right)^{c}$, we have
- if $f(x)=0$, then $h(x)=0$;
. if $f(x)>0$, then $h(x)>0$;
- if $f(x)<0$, then $h(x)<0$.

Since $Z\left(s_{1}\right)^{c} \cap Z\left(s_{2}\right)^{c} \cap Z\left(s_{3}\right)^{c}$ is a dense subset, by above claim, $\bar{h}=\bar{f}$.

- Let $f, g \in A$ and $\overline{h_{1}}, \overline{h_{2}} \in \bar{f}+\bar{g}$ in $A /{ }_{m} T$. By an argument similar of the preceding item, the signals of $h_{1}, h_{2}$ are equal in dense subset and thus $\overline{h_{1}}=\overline{h_{2}}$.


## Chapter 3

## From Multirings to Superrings

The concept of superring first appears in ([6]). The very first advantage of considering superrings instead of hyperrings is the possibility of built a theory of polynomials and matrices, available for hyperrings but only closed by constructions in superrings. There are many important advances and results in superring theory, and for instance, we recommend for example, the following papers: [3], [5], [6, [4], 49], 54], 51, 50].

Surprisingly we have obtained an interesting theory of matrices, linear systems, vector spaces and algebraic extensions available for a certain subclass of superfields. If $R$ is a full superring, then $M_{m \times n}(R)$ and $R[X]$ are superrings (Theorem 3.2.6 and 3.4.2). We also obtained a kind of simple algebraic extension for a superfield $F$ (Theorem 3.6.12), which culminate in the existence and unicity of a full algebraic extension of a superfield $F$ (Theorems 3.7.3 and 3.7.4. If $F$ is a linearly closed superfield (the system $A x=0$ always have a non trivial solution), then we have a well defined dimension theory for the vector spaces over $F$ (Theorem 3.8.21). The main examples of linearly closed superfields are hyperbolic hyperfields (3.8.23) and simple full algebraic extensions over a linearly closed superfield (3.8.25). The linearly closed interpreted in the context of special groups leads to interesting Isotropic (Corollary 3.8.27) and Hyperbolic (Corollary 3.8.28) interpolations.

We finish this Chapter with a quantifier elimination procedure for superfields (Theorem 3.9.3), which is a direct generalization of a result obtained in [19.

### 3.1 Superrings, Superfields

Definition 3.1.1 (Definition 5 in [6]). An associative superring is a structure $(S,+, \cdot,-, 0,1)$ such that:
$i-(S,+,-, 0)$ is a commutative multigroup.
ii - $(S, \cdot, 1)$ is a multimonoid.
iii - 0 is an absorbing element: $a \cdot 0=\{0\}=0 \cdot a$, for all $a \in S$.
iv - The weak/semi distributive law holds: if $d \in c .(a+b)$ and $e \in(a+b) c$ then $d \in c a+c b$ and $e \in a c+b c$, for all $a, b, c, d, e \in S$.
$v$ - The rule of signals holds: $-(a b)=(-a) b=a(-b)$, for all $a, b \in S$.
A superdomain is a non-trivial superring without zero-divisors in this new context, i.e. whenever

$$
0 \in a \cdot b \text { iff } a=0 \text { or } b=0
$$

A quasi-superfield is a non-trivial superring such that every nonzero element is invertible in this new contex ${ }^{1}$, i.e. whenever

$$
\text { For all } a \neq 0 \text { exists } b \text { such that } 1 \in a \cdot b \text {. }
$$

A superfield is a quasi-superfield which is also a superdomain. A superring is full if for all $a, b, c, d \in$ $S, d \in c \cdot(a+b)$ iff $d \in c a+c b$.

A superring, superdomain or superfield is commutative (associative) if $(S, \cdot, 1)$ is respectively commutative (associative).

From now on, all superrings will be commutative (with exceptions sinalized).

Example 3.1.2. Every multiring can be seen as a superring, in the very same fashion of 1.2.9(a). Our main example of superring is the superring of multipolynomials $R[X]$ over a multiring $R$. The construction will be presented in short in Section 3.4. For more details, see [15], [6] or [11].

Now we treat about morphisms.

Definition 3.1.3. Let $A$ and $B$ superrings. $A \operatorname{map} f: A \rightarrow B$ is a morphism if for all $a, b, c \in A$ :

$$
\begin{array}{ll}
i-f(0)=0 ; & i v-c \in a+b \Rightarrow f(c) \in f(a)+f(b) ; \\
i i-f(1)=1 ; & v-c \in a \cdot b \Rightarrow f(c) \in f(a) \cdot f(b)
\end{array}
$$

$A$ morphism $f$ is a full morphism if for all $a, b \in A$,

$$
f(a+b)=f(a)+f(b) \text { and } f(a \cdot b)=f(a)+f(b)
$$

From now on, we use the following conventions: Let $(R,+, \cdot,-, 0,1)$ be a superring, $p \in \mathbb{N}$ and consider a $p$-tuple $\vec{a}=\left(a_{0}, a_{1}, \ldots, a_{p-1}\right)$. We define the finite sum by:

$$
\begin{aligned}
& x \in \sum_{i<0} a_{i} \text { iff } x=0 \\
& x \in \sum_{i<p} a_{i} \text { iff } x \in y+a_{p-1} \text { for some } y \in \sum_{i<p-1} a_{i} \text {, if } p \geq 1 .
\end{aligned}
$$

and the finite product by:

$$
\begin{aligned}
& x \in \prod_{i<0} a_{i} \text { iff } x=1 \\
& x \in \prod_{i<p} a_{i} \text { iff } x \in y \cdot a_{p-1} \text { for some } y \in \prod_{i<p-1} a_{i}, \text { if } p \geq 1
\end{aligned}
$$

Thus, if $\left(\vec{a}_{0}, \vec{a}_{1}, \ldots, \vec{a}_{p-1}\right)$ is a $p$-tuple of tuples $\vec{a}_{i}=\left(a_{i 0}, a_{i 1}, \ldots, a_{i m_{i}}\right)$, then we have the finite

[^6]sum of finite products:
\[

$$
\begin{aligned}
& x \in \sum_{i<0} \prod_{j<m_{i}} a_{i j} \text { iff } x=0, \\
& x \in \sum_{i<p} \prod_{j<m_{i}} a_{i j} \text { iff } x \in y+z \text { for some } y \in \sum_{i<p-1} \prod_{j<m_{i}} a_{i j} \text { and } z \in \prod_{j<m_{p-1}} a_{p-1, j}, \text { if } p \geq 1 .
\end{aligned}
$$
\]

Remark 3.1.4. Note that, in this sense, we have (for example) that

$$
\prod_{j=1}^{4} a_{j}=\left(\left(a_{1} a_{2}\right) a_{3}\right) a_{4}
$$

lembasic1
Lemma 3.1.5 (Basic Facts). Let $A$ be a superring.
$a-$ For all $n \in \mathbb{N}$ and all $a_{1}, \ldots, a_{n} \in A$, the sum $a_{1}+\ldots+a_{n}$ does not depends on the order of the entries, and if is $A$ is associative, the product $a_{1} \cdot \ldots \cdot a_{n}$ also does not depends on the order of the entries.
$b$ - If $A$ is a full superdomain, then $a x=a y$ for some $a \neq 0$ imply $x=y$.
$c$ - If $A$ is full, then for all $d, a_{1}, \ldots, a_{n} \in A$

$$
d\left(a_{1}+\ldots+a_{n}\right)=d a_{1}+\ldots+d a_{n} .
$$

$d$ - Suppose $A$ is a full superdomain and let $a \in A \backslash\{0\}$. If $1 \in(a \cdot b) \cap(a \cdot c)$ then $b=c$.
$e-(N e w t o n ' s ~ B i n o m ~ F o r m u l a) ~ F o r ~ n \geq 1 ~ a n d ~ X \subseteq A ~ d e n o t e ~$

$$
n X:=\sum_{i=1}^{n} X .
$$

Then for $A, B \subseteq A$,

$$
(A+B)^{n} \subseteq \sum_{j=0}^{n}\binom{n}{j} A^{j} B^{n-j}
$$

Proof.
a - It is an immediate consequence of associativity and induction.
b - Let $a x=a y$ for some $a \neq 0$. Then $a x-a y=a y-a y$. Since $A$ is full, $a(x-y)=a y-a y$, and then,

$$
0 \in a y-a y=a(x-y) .
$$

Moreover, $0 \in a z$ for some $z \in x-y$. Since $A$ is a superdomain and $a \neq 0, z=0$. Then $0 \in x-y$, which imply $x=y$.
c - By induction, we only need to proof the case $n=2$. Let $a, b, c, d \in A$. We already know that $d(a+b+c) \subseteq d a+d b+d c$. Now consider $x \in d a+d b+d c$. Then $x \in e+d c$ for some $e \in d a+d b=d(a+b)$. Then $e \in d e^{\prime}$ with $e^{\prime} \in a+b$ and $x \in e+d c \subseteq d e^{\prime}+d c=d\left(e^{\prime}+c\right)$. Hence

$$
x \in d\left(e^{\prime}+c\right) \subseteq d(a+b+c) .
$$

d - Let $1 \in(a \cdot b) \cap(a \cdot c)$. Then

$$
0 \in 1-1 \subseteq(a \cdot b)-(a \cdot c)=a \cdot(b-c)
$$

Since $0 \in a \cdot(b-c)$ and $a \neq 0$ we have $0 \in b-c$, which imply $b=c$.
$\mathrm{e}-\mathrm{By}$ induction is enough to prove the case $n=2$. We have

$$
\begin{aligned}
(A+B)^{2} & :=(A+B)(A+B) \subseteq A(A+B)+B(A+B) \subseteq A^{2}+A B+B A+B^{2} \\
& =A^{2}+A B+A B+B^{2}=A^{2}+2 A B+B^{2}:=\sum_{j=0}^{2}\binom{n}{j} A^{j} B^{n-j}
\end{aligned}
$$

factstrong 2
Lemma 3.1.6 (Facts about full morphisms of superrings). Let $f: A \rightarrow B$ be a full morphism of superrings. Then
$a-$ For all $a_{1}, \ldots, a_{n} \in A$,

$$
f\left(a_{1}+\ldots+a_{n}\right)=f\left(a_{1}\right)+\ldots+f\left(a_{n}\right) .
$$

$b-$ For all $a_{1}, \ldots, a_{n}, b_{1}, \ldots, b_{n} \in A$,

$$
f\left[\left(a_{1}+b_{1}\right)\left(a_{2}+b_{2}\right) \ldots\left(a_{n}+b_{n}\right)\right]=\left(f\left(a_{1}\right)+f\left(b_{1}\right)\right)\left(f\left(a_{2}\right)+f\left(b_{2}\right)\right) \ldots\left(f\left(a_{n}\right)+f\left(b_{n}\right)\right)
$$

$c-$ For all $c_{1}, \ldots, c_{n}, d_{1}, \ldots, d_{n} \in A$,

$$
f\left(c_{1} d_{1}+c_{2} d_{2}+\ldots+c_{n} d_{n}\right)=f\left(c_{1}\right) f\left(d_{1}\right)+f\left(c_{2}\right) f\left(d_{2}\right)+\ldots+f\left(c_{n}\right) f\left(d_{n}\right)
$$

$d$ - For all $a_{0}, \ldots, a_{n}, \alpha \in A$,

$$
f\left(a_{0}+a_{1} \alpha+\ldots+a_{n} \alpha^{n}\right)=f\left(a_{0}\right)+f\left(a_{1}\right) f(\alpha)+\ldots+f\left(a_{n}\right) f(\alpha)^{n} .
$$

Here we always interpret $a b^{n}:=a\left(b^{n}\right)$, unless stated contrary.
$e-L e t A_{1}, A_{2}, A_{3}$ be superrings with injective morphisms (embeddings) $i_{12}: A_{1} \rightarrow A_{2}, i_{13}: A_{1} \rightarrow$ $A_{3}$ and $i_{23}: A_{2} \rightarrow A_{3}$.


Suppose that $i_{13}=i_{23} \circ i_{12}$ is a full embedding. If $i_{23}$ is a full embedding then $i_{12}$ is a full embedding.

## Definition 3.1.7.

$i$ - The characteristic of a superring is the smaller integer $n \geq 1$ such that

$$
0 \in \sum_{i<n} 1,
$$

otherwise the characteristic is zero. For full superdomains, this is equivalent to say that $n$ is the smaller integer such that

$$
\text { For all } a, 0 \in \sum_{i<n} a \text {. }
$$

ii - An ideal of a superring $A$ is a non-empty subset $\mathfrak{a}$ of $A$ such that $\mathfrak{a}+\mathfrak{a} \subseteq \mathfrak{a}$ and $A \mathfrak{a} \subseteq \mathfrak{a}$. We denote

$$
\Im(A)=\{I \subseteq A: I \text { is an ideal }\} .
$$

iii - Let $S$ be a subset of a superring $A$. We define the ideal generated by $S$ as

$$
\langle S\rangle:=\bigcap\{\mathfrak{a} \subseteq A \text { ideal }: S \subseteq \mathfrak{a}\} .
$$

If $S=\left\{a_{1}, \ldots, a_{n}\right\}$, we easily check that

$$
\left\langle a_{1}, \ldots, a_{n}\right\rangle=\sum A a_{1}+\ldots+\sum A a_{n}, \text { where } \sum A a=\bigcup_{n \geq 1}\{\underbrace{A a+\ldots+A a}_{n \text { times }}\} .
$$

Note that if $A$ is a full superring, then $\sum A a=A a$.
iv - An ideal $\mathfrak{p}$ of $A$ is said to be prime if $1 \notin \mathfrak{p}$ and $a b \subseteq \mathfrak{p} \Rightarrow a \in \mathfrak{p}$ or $b \in \mathfrak{p}$. We denote

$$
\operatorname{Spec}(A)=\{\mathfrak{p} \subseteq A: \mathfrak{p} \text { is a prime ideal }\} .
$$

$v-$ An ideal $\mathfrak{p}$ of $A$ is said to be strongly prime if $1 \notin \mathfrak{p}$ and $a b \cap \mathfrak{p} \neq \emptyset \Rightarrow a \in \mathfrak{p}$ or $b \in \mathfrak{p}$. We denote

$$
\operatorname{Spec}_{s}(A)=\{\mathfrak{p} \subseteq A: \mathfrak{p} \text { is a strongly prime ideal }\} .
$$

Note that every strongly prime ideal is prime.
$v i$ - An ideal $\mathfrak{m}$ is maximal if it is proper and for all ideals $\mathfrak{a}$ with $\mathfrak{m} \subseteq \mathfrak{a} \subseteq A$ then $\mathfrak{a}=\mathfrak{m}$ or $\mathfrak{a}=A$.
vii - For an ideal $I \subseteq A$, we define operations in the quotient $A / I=\{x+I: x \in A\}=\{\bar{x}: x \in A\}$, by the rules

$$
\begin{aligned}
\bar{x}+\bar{y} & =\{\bar{z}: z \in x+y\} \\
\bar{x} \cdot \bar{y} & =\{\bar{z}: z \in x y\}
\end{aligned}
$$

for all $\bar{x}, \bar{y} \in A / I$.

## Remark 3.1.8.

$a-$ If $A$ is a multiring, then every prime ideal is strongly prime. We do not know if this is the case for general superrings.
$b$ - If $A$ is a multiring, then every maximal ideal is prime (Proposition 1.7 of [23]). For a general superring $A$, we do not know if a maximal ideal is prime.
c - In his Ph.D Thesis [23], H. Ribeiro deals with elements weakly invertible on a multiring A. This could be an alternative in dealing with the above questions.
With all conventions and notations above, we obtain the following Lemma, which recover for superrings some properties holding for rings (and multirings).

Lemma 3.1.9. Let $A$ be an associative superring and $I$ an ideal.
$i-I=A$ if and only if $1 \in I$.
ii $-A / I$ is a superring. Moreover, if $A$ is full then $A / I$ is also full.
iii - $I$ is strongly prime if and only if $A / I$ is a superdomain.
If $A$ is full, then
iv $-I=A$ if and only if $1 \in I$, which occurs if and only if $A^{*} \cap I \neq \emptyset$ (in other words, if and only if I contains an invertible element).
$v-A$ is a superfield if and only if $\mathfrak{I}(A)=\{0, A\}$.
vi - $I$ is maximal if and only if $A / I$ is a superfield.
Proposition 3.1.10. Let $A$ be an associative superring and $I$ an ideal.
$i$ - If $I$ is a maximal ideal, then it is prime.
ii - The ideal $I$ is prime if, and only if, $A / I$ is quasi-superdomain ${ }^{2}$.
iii - (Prime Ideal Theorem) Let $S \subseteq A$ be a multiplicative set $(1 \in S$ and $S \cdot S \subseteq S$ ). Suppose that $S \cap I=\emptyset$. Then there is a prime ideal $p$ such that $I \subseteq p$ and $S \cap p=\emptyset$.

Proof.
i - Let $a, b \in A$ with $a b \subseteq I$. Assume that $a \notin I$ and consider the ideal $J=I+(a)$. Then there are $x \in I, t_{1}, \ldots, t_{n} \in A$ such that $1 \in x+a t_{1}+\cdots+a t_{n}$. Thus $b \in b x+b a t_{1}+\cdots+b a t_{n} \subseteq I$.
ii - If $I$ is prime and $\bar{a} \cdot \bar{b}=\{\overline{0}\}$ in $A / I$, then $a \cdot b \subseteq I$. Therefore, by primality, $a \in I$ or $b \in I$. Thus $\bar{a}=0$ or $\bar{b}=0$ in $A / I$. Reciprocally, assume $A / I$ a quasi-superdomain and let $a, b \in A$ with $a b \subseteq I$. Then $\bar{a} \cdot \bar{b}=\{\overline{0}\}$ and by hypothesis follows $a \in I$ or $b \in I$, as desired.
iii - Consider the partial order $\mathcal{X}=\{J \subseteq A: J$ is an ideal and $S \cap J=\emptyset\}$ ordered by inclusion. Since the directed union of ideals is again an ideal, we have by Zorn's Lemma that $\mathcal{X}$ has a maximal element $p$. Suppose that $p$ is not prime, that it, there is $a, b \in$ with $a b \subseteq p$ and $a, b \notin p$. Now notice that $J_{a}=p+(a)$ and $J_{b}=p+(b)$ are ideals that proper extend $p$. Hence, by maximality, there are $s, v \in S, x, y \in p, t_{1}, \ldots, t_{n}, w_{1}, \ldots, w_{k} \in A$ with

$$
\begin{aligned}
& s \in x+a t_{1}+\cdots+a t_{n} \\
& v \in y+b w_{1}+\cdots+b w_{k}
\end{aligned}
$$

These equations implies that

$$
s v \subseteq x y+x b w_{1}+\cdots+x b w_{k}+y a t_{1}+\cdots+y a t_{n}+\sum_{i, j} a b t_{i} w_{j} \subseteq S \cap p
$$

a contradiction. Then $p$ is prime.

[^7]
### 3.2 Matrices and determinants over commutative superrings

Definition 3.2.1. Let $m, n$ be positive integers. A $m \times n$ matrix over a commutative superring $R$ is a double sequence $A$ of elements of $F$, distributed in $m$ rows and $n$ columns. The set of $m \times n$ matrices is denoted by $M_{m \times n}(R)$. When $m=n$, we denote $n \times n$ matrices by $M_{n}(R)$.

A matrix $A \in M_{m \times n}(R)$ is represented simply by $A=\left(a_{i j}\right)$ (with $m$ and $n$ subscript if necessary) or by a table as below:

$$
A=\left(\begin{array}{cccc}
a_{11} & a_{12} & \ldots & a_{1 n} \\
a_{21} & a_{22} & \ldots & a_{2 n} \\
\vdots & \vdots & \ddots & \vdots \\
a_{m 1} & a_{m 2} & \ldots & a_{m n}
\end{array}\right)
$$

For $A, B \in M_{n \times m}(R)$ and $\lambda \in R$ with $A=\left(a_{i j}\right)$ and $B=\left(b_{i j}\right)$ we define $-A=\left(-a_{i j}\right)$ and (multi) operations

$$
A+B:=\left\{\left(d_{i j}\right): d_{i j} \in a_{i j}+b_{i j} \text { for all } i, j\right\} \neq \emptyset
$$

and

$$
\lambda A=\left\{\left(d_{i j}\right): d_{i j} \in \lambda a_{i j} \text { for all } i, j\right\}
$$

If $A \in M_{n \times m}(R)$ with $A=\left(a_{i j}\right)$ and $B \in M_{m \times p}(R)$ with $B=\left(a_{i j}\right)$, we define

$$
A \cdot B=A B \subseteq M_{n \times p}(R)
$$

by

$$
A B=\left\{\left(d_{i j}\right): d_{i j} \in \sum_{k=1}^{n} a_{i k} b_{k j}=a_{i 1} b_{1 j}+a_{i 2} b_{2 j}+\ldots+a_{i n} b_{n k} \text { for all } i, j\right\} \neq \emptyset
$$

We denote $0=\left(0_{i j}\right) \in M_{m \times n}(R)$ and $1=\left(\delta_{i j}\right)_{i j} \in M_{n}(R)$ the usual zero and identity matrices respectively.

We say that $A \in M_{n}(R)$ is invertible iff there exist $B \in M_{n}(R)$ with $1 \in A B$ and $1 \in B A$.
Of course, we adopt freely the usual simplified notation from commutative algebra. For example for $A=\left(a_{i j}\right)$ and $B=\left(b_{i j}\right)$ we simply write

$$
\left(\begin{array}{cccc}
a_{11} & a_{12} & \ldots & a_{1 n} \\
a_{21} & a_{22} & \ldots & a_{2 n} \\
\vdots & \vdots & \ddots & \vdots \\
a_{m 1} & a_{m 2} & \ldots & a_{m n}
\end{array}\right)+\left(\begin{array}{cccc}
b_{11} & b_{12} & \ldots & b_{1 n} \\
b_{21} & b_{22} & \ldots & b_{2 n} \\
\vdots & \vdots & \ddots & \vdots \\
b_{m 1} & b_{m 2} & \ldots & b_{m n}
\end{array}\right)=\left(\begin{array}{cccc}
a_{11}+b_{11} & a_{12}+b_{12} & \ldots & a_{1 n}+b_{1 n} \\
a_{21}+b_{21} & a_{22}+b_{22} & \ldots & a_{2 n}+b_{2 n} \\
\vdots & \vdots & \ddots & \vdots \\
a_{m 1}+b_{m 1} & a_{m 2}+b_{m 2} & \ldots & a_{m n}+b_{m n}
\end{array}\right)
$$

with the analogous simplifications for $\lambda A$ and $A B$.

Example 3.2.2. Consider $X_{2}=\{-2,-1,0,1,2\}$ as in Example 1.2.12 and matrices $A, B, C \in$ $M_{2}\left(X_{2}\right)$ given by

$$
A=\left(\begin{array}{ll}
1 & 1 \\
0 & 1
\end{array}\right), B=\left(\begin{array}{cc}
-1 & 1 \\
0 & -1
\end{array}\right) \text { and } C=\left(\begin{array}{cc}
2 & 0 \\
-1 & 2
\end{array}\right)
$$

With our notations we have

$$
\left.\left.\begin{array}{l}
A+B=\left(\begin{array}{ll}
1-1 & 1+1 \\
0+0 & 1-1
\end{array}\right)= \\
\left\{\left(\begin{array}{cc}
-1 & 1 \\
0 & -1
\end{array}\right),\left(\begin{array}{cc}
-1 & 1 \\
0 & 0
\end{array}\right),\left(\begin{array}{cc}
-1 & 1 \\
0 & 1
\end{array}\right),\left(\begin{array}{cc}
0 & 1 \\
0 & -1
\end{array}\right),\left(\begin{array}{ll}
0 & 1 \\
0 & 0
\end{array}\right),\left(\begin{array}{ll}
0 & 1 \\
0 & 1
\end{array}\right),\left(\begin{array}{cc}
1 & 1 \\
0 & -1
\end{array}\right),\left(\begin{array}{ll}
1 & 1 \\
0 & 0
\end{array}\right),\left(\begin{array}{ll}
1 & 1 \\
0 & 1
\end{array}\right)\right\} \\
A \cdot B=\binom{1 \cdot(-1)+1 \cdot 0}{0 \cdot(-1)+1 \cdot(0)} 0 \cdot 1+1 \cdot(-1) \\
0
\end{array}\right)=\left(\begin{array}{cc}
-1 & 1-1 \\
0 & -1
\end{array}\right)=,\left(\begin{array}{cc}
-1 & -1 \\
0 & -1
\end{array}\right),\left(\begin{array}{cc}
-1 & 0 \\
0 & -1
\end{array}\right),\left(\begin{array}{cc}
-1 & 1 \\
0 & -1
\end{array}\right)\right\} \text {. }
$$

Here we will "export" the usual terminology of diagonal, triangular, block and elementary matrices available for fields to superfields.

With these, using the fact that $(R,+,-, 0)$ is a commutative multigroup we immediately have the following Theorem.
matrix1
Theorem 3.2.3. Let $R$ be a superring. Then $\left(M_{m \times n}(R),+,-, 0\right)$ is a commutative multigroup.
For a general associative superring $R$, the matrix product in $M_{n}(R)$ is not associative in general (and of course, $M_{n}(R)$ is not an associative superring in general).

Example 3.2.4. Let $R=X_{2}$ as in Example 1.2.12. Of course, $R$ is not full because, for example,

$$
2(1-1)=\{-2,0,2\} \text { and } 2-2=R .
$$

Now, consider the matrices

$$
A=\left(\begin{array}{ll}
1 & 1 \\
0 & 1
\end{array}\right), B=\left(\begin{array}{cc}
-1 & 1 \\
0 & -1
\end{array}\right) \text { and } C=\left(\begin{array}{cc}
2 & 0 \\
-1 & 2
\end{array}\right) .
$$

In fact we have

$$
\begin{aligned}
(A B) C & =\left[\left(\begin{array}{ll}
1 & 1 \\
0 & 1
\end{array}\right)\left(\begin{array}{cc}
-1 & 1 \\
0 & -1
\end{array}\right)\right]\left(\begin{array}{cc}
2 & 0 \\
-1 & 2
\end{array}\right) \\
& =\left(\begin{array}{cc}
-1 & 1-1 \\
0 & -1
\end{array}\right)\left(\begin{array}{cc}
2 & 0 \\
-1 & 2
\end{array}\right)=\left(\begin{array}{cc}
-2-(1-1) & 2(1-1) \\
1 & -2
\end{array}\right)
\end{aligned}
$$

and

$$
\begin{aligned}
A(B C) & =\left(\begin{array}{ll}
1 & 1 \\
0 & 1
\end{array}\right)\left[\left(\begin{array}{cc}
-1 & 1 \\
0 & -1
\end{array}\right)\left(\begin{array}{cc}
2 & 0 \\
-1 & 2
\end{array}\right)\right] \\
& =\left(\begin{array}{ll}
1 & 1 \\
0 & 1
\end{array}\right)\left(\begin{array}{cc}
-2-1 & 2 \\
1 & -2
\end{array}\right)=\left(\begin{array}{cc}
(-2-1)+1 & 2-2 \\
1 & -2
\end{array}\right) .
\end{aligned}
$$

Then we have $(A B) C \neq A(B C)$.

Despite the fact that $M_{n}(R)$ is not an associative superring in general, it is a structure of interest (as we will see, for example, in Theorem 3.6.12) and in the following Lemma we collect the properties holding for the product in the general case.

Lemma 3.2.5. Let $R$ be a superring and $A, B, C$ matrices in $M_{n}(R)$. Then:
$a-A \cdot 0=0 \cdot A=\{0\}$.
$b$ - If $m=n$, then $A \cdot 1=1 \cdot A=\{A\}$.
$c$ - If $A(B+C) \subseteq A B+A C$, with equality if $R$ is full.
$d-(B+C) A \subseteq B A+C A$, with equality if $R$ is full.
$e$ - If $R$ is associative and full, then $(A B) C=A(B C)$.
Firstly, let us explicit the notation for $A B$ : we write

$$
A B=\left(\begin{array}{cccc}
D_{11} & D_{12} & \ldots & D_{1 n} \\
D_{21} & D_{22} & \ldots & D_{2 n} \\
\vdots & \vdots & \ddots & \vdots \\
D_{m 1} & D_{m 2} & \ldots & D_{m n}
\end{array}\right)
$$

or $A B=\left(D_{i j}\right)$ where for all $i, j, D_{i j}$ is the set

$$
D_{i j}=\sum_{k=1}^{n} a_{i k} b_{k j}=a_{i 1} b_{1 j}+a_{i 2} b_{2 j}+\ldots+a_{i n} b_{n k} \text { (of course, this is an equality of sets). }
$$

Alternatively, we can proceed more directly, simply writing

$$
A B=\left(\begin{array}{cccc}
\sum_{k=1}^{n} a_{1 k} b_{k 1} & \sum_{k=1}^{n} a_{1 k} b_{k 2} & \ldots & \sum_{k=1}^{n} a_{1 k} b_{k n} \\
\sum_{k=1}^{n=1} a_{2 k} b_{k 1} & \sum_{k=1}^{n} a_{2 k} b_{k 2} & \ldots & \sum_{k=1}^{n} a_{2 k} b_{k n} \\
\vdots & \vdots & \ddots & \vdots \\
\sum_{k=1}^{n} a_{m k} b_{k 1} & \sum_{k=1}^{n} a_{m k} b_{k 2} & \ldots & \sum_{k=1}^{n} a_{m k} b_{k n}
\end{array}\right)
$$

Now we proceed with the proof of the Lemma.
Proof of Lemma 3.2.5. The argument here is in fact very similar to those one used in linear algebra over fields.
a - Let $A \cdot 0=\left(D_{i j}\right)$ (as explained above). For all $i, j$ we have

$$
D_{i j}=\sum_{k} a_{i k} 0_{k j}=\sum_{k} a_{i k} \cdot 0=\sum_{k} 0=0 .
$$

Then $A \cdot 0=\{0\}$. The same reasoning provide $0 \cdot A=\{0\}$.
b- Let $A \cdot 1=\left(D_{i, j}\right)$. Since $1=\left(\delta_{i, j}\right)$ with $\delta_{i, j} \in\{0,1\}$ and $\delta_{i, j}=1$ iff $i=j$, we have

$$
D_{i, j}=\sum_{k} a_{i, k} \delta_{k, j}=a_{i, j} \cdot 1=\left\{a_{i, j}\right\} .
$$

Thus, $A \cdot 1=\{A\}$. Similarly, $1 \cdot A=\{A\}$.
c -

$$
\begin{aligned}
A(B+C) & =\left(a_{i, j}\right)_{i, j} \cdot\left(b_{i, j}+c_{i, j}\right)_{i, j} \\
& =\left(\sum_{k} a_{i, k}\left(b_{k, j}+c_{k, j}\right)\right)_{i, j} \\
& \subseteq\left(\sum_{k} a_{i, k} b_{k, j}+a_{i, k} c_{k, j}\right)_{i, j} \\
& =\left(\sum_{k} a_{i, k} b_{k, j}\right)_{i, j}+\left(\sum_{k} a_{i, k} c_{k, j}\right)_{i, j}=A B+A C .
\end{aligned}
$$

When $R$ is full, it is immediate from above that $A(B+C)=A B+A C$.
d - Similar argument as above.
e - Assume that $R$ is associative and full. Then

$$
\begin{aligned}
(A B) C & =\left[\left(a_{i, j}\right)_{i, j} \cdot\left(b_{i, j}\right)_{i, j}\right] \cdot\left(c_{i, j}\right)_{i, j} \\
& =\left(\sum_{k} a_{i, k} b_{k, j}\right)_{i, j} \cdot\left(c_{i, j}\right)_{i, j} \\
& =\left(\sum_{l}\left(\sum_{k} a_{i, k} b_{k, l}\right) \cdot c_{l, j}\right)_{i, j} \\
& =\left(\sum_{l} \sum_{k} a_{i, k} b_{k, l} c_{l, j}\right)_{i, j} \\
& =\left(\sum_{k} a_{i, k} \cdot\left(\sum_{l} b_{k, l} c_{l, j}\right)\right)_{i, j} \\
& =\left(a_{i, j}\right) \cdot\left(\sum_{l} b_{i, l} c_{l, j}\right)_{i, j}=A(B C) .
\end{aligned}
$$

In fact, with Theorem 3.2.3 and Lemma 3.2.5 we conclude the following.
Theorem 3.2.6. For a superring $R$, if $R$ is associative and full then $M_{n}(R)$ is a full superring, that is non-commutative if $n \geq 2$.

We also have a generalized version of Lemma 3.2 .5 (with the proof similar to the one given there).
matrix4
Lemma 3.2.7. Let $R$ be a superring and $A, B, C, D, E, F$ matrices with $A \in M_{m \times n}(R), B, C \in$ $M_{n \times p}(R), D, E \in M_{p \times m}(R)$ and $F \in M_{p \times q}(R)$. Then:
$a-A \cdot 0_{n \times p}=\left\{0_{m \times p}\right\}$ and $0_{p \times n} \cdot A=\left\{0_{p \times n}\right\}$.
$b-A \cdot 1_{n \times n}=1_{m \times m} \cdot A=\{A\}$.
$c-A(B+C) \subseteq A B+A C$, with equality if $R$ is full.
$d-(D+E) A \subseteq D A+E A$, with equality if $R$ is full.
$e$ - If $R$ is associative and full, then $(A B) F=A(B F)$.
Despite the fact that full associativity do not hold in $M_{n}(R)$ for a general superring $R$, we have the following useful results. We start with a technical Definition:

Definition 3.2.8. Let $R$ be a superring. We say that $R$ is proto-full if for all $a, b, c, d \in R$

$$
[(a b+a c) d] \cap[a(b d+c d)] \neq \emptyset .
$$

Lemma 3.2.9. Let $R$ be an associative and proto-full superring. Then for all $a, b_{1}, \ldots, b_{n}, d \in R$ we have

$$
\left[\left(a b_{1}+\ldots+a b_{n}\right) d\right] \cap\left[a\left(b_{1} d+\ldots+b_{n} d\right)\right] \neq \emptyset .
$$

Then rewriting the proof of Lemma $3 \cdot 2.7$ (e) we get the following.
matrix5
Lemma 3.2.10. Let $R$ be an associative and proto full superring and $A, B, C$ matrices with $A \in$ $M_{m \times n}(R), B \in M_{n \times p}(R)$ and $C \in M_{p \times q}(R)$. Then

$$
[(A B) C] \cap[A(B C)] \neq \emptyset
$$

determinant
Definition 3.2.11. Let $\mathcal{A} \subseteq M_{n}(R)$. We define the determinant of $\mathcal{A}$ as the subset $\operatorname{det}(\mathcal{A}) \subseteq R$ given by the rule

$$
\operatorname{det}(\mathcal{A})=\bigcup_{A \in \mathcal{A}}\left\{\sum_{\sigma \in S_{n}} \operatorname{sgn}(\sigma) \prod_{j=1}^{n} a_{j \sigma(j)}\right\} .
$$

If $\mathcal{A}=\{A\}$ we simply write $\operatorname{det}(A)$ for the above formula.
Lemma 3.2.12 (Properties of Determinant). Let $R$ be associative and $A, B \in M_{n}(R), A=\left(a_{i j}\right)$, $B=\left(b_{i j}\right)$ and $\lambda \in R$. Then:
$a-\operatorname{det}(\lambda A) \subseteq \lambda^{n} \operatorname{det}(A)$, with equality if $R$ is full;
$b$ - if there is an entire row or column of zeros in $A$ then $\operatorname{det}(A)=\{0\}$.
$c-$ if $A=\left(a_{i j}\right)$ is a triangular matrix (and in particular, diagonal matrix) then $\operatorname{det}(A)=a_{11} a_{22} \ldots a_{n n}$.
Proof.
a - Using the very Definition we get

$$
\begin{aligned}
\operatorname{det}(\lambda A) & =\sum_{\sigma \in S_{n}} \operatorname{sgn}(\sigma) \prod_{j=1}^{n}\left(\lambda a_{j \sigma(j)}\right)=\sum_{\sigma \in S_{n}} \operatorname{sgn}(\sigma) \lambda^{n} \prod_{j=1}^{n} a_{j \sigma(j)} \\
& \subseteq \lambda^{n}\left(\sum_{\sigma \in S_{n}} \operatorname{sgn}(\sigma) \prod_{j=1}^{n} a_{j \sigma(j)}\right)=\lambda^{n} \operatorname{det}(A)
\end{aligned}
$$

b-In this case we have $0 \in\left\{a_{1 \sigma(1)}, a_{2 \sigma(2)}, \ldots, a_{n \sigma(n)}\right\}$ for all $n \in S_{n}$. Then

$$
\prod_{j=1}^{n} a_{j \sigma(j)}=\{0\} \text { for all } \sigma \in S_{n}
$$

implying $\operatorname{det}(A)=\{0\}$.
c - Follow immediately from Definition 3.2.11.

### 3.3 Linear systems over superfields

Throughout this entire Section fix a superfield $F$.
lin-equation
Definition 3.3.1. A linear equation is an equation (term in the language of superfields) of type

$$
A x \subseteq B
$$

where $A \in M_{1 \times n}(F)(n \in \mathbb{N})$, $x$ is a $n \times 1$ vector of variables and $B \subseteq F$.
Remark 3.3.2. We are defining "linear equations" (and more lately, "linear systems") in terms of matrices. We do this because for superfields we cannot agglutinate scalars and variables as we do in general linear algebra. For example, there is no reason for " $a_{1} x_{1}+b_{1} x_{1}$ " be equal to " $\left(a_{1}+b_{1}\right) x_{1}, 3$,

Despite this Remark, given a linear equation $A x \subseteq B$, we can "colloquially" write

$$
a_{1} x_{1}+a_{2} x_{2}+\ldots+a_{n} x_{n} \subseteq B,
$$

with $a_{1}, \ldots, a_{n} \in F$. Of course, we could consider the equation

$$
a_{1} x_{1}+a_{2} x_{2}+\ldots+a_{n} x_{n} \supseteq B
$$

as a linear one and proceed with two types of linear equations. But the type considered in 3.3.1 seems to be (at first sight) more "natural".

We can use to this "coloquial" to write our equations (and further, systems), and while we are dealing with one equation (system), we will proceed with this "coloquial" language. But in order to get more general proofs and Definitions, we will always proceed with matrices.

Definition 3.3.3. A solution (weak solution) of a linear equation $A x \subseteq B$, is a matrix $d \in$ $M_{n \times 1}(F)$ such that $A d \subseteq B \quad(A d \cap B \neq \emptyset)$.

Definition 3.3.4. A linear system is a conjunction of equations (term in the language of superfields) of type

$$
A x \subseteq B
$$

where $A \in M_{m \times n}(F)(n \in \mathbb{N})$, $x$ is a $n \times 1$ vector of variables and $B \subseteq M_{m \times 1}(F)$.

[^8]In this sense, a Linear system can be colloquially represented as usual

$$
S:\left\{\begin{array}{l}
a_{11} x_{1}+a_{12} x_{2}+\ldots+a_{1 n} x_{n} \subseteq B_{1} \\
a_{21} x_{1}+a_{22} x_{2}+\ldots+a_{2 n} x_{n} \subseteq B_{2} \\
\vdots \\
a_{m 1} x_{1}+a_{m 2} x_{2}+\ldots+a_{m n} x_{n} \subseteq B_{m}
\end{array}\right.
$$

A (weak) solution of a Linear system is a tuple $\left(d_{1}, \ldots, d_{n}\right)$ such that $A d \subseteq B(A d \cap B \neq \emptyset)$.
Definition 3.3.5. A Linear system $A x \subseteq B$ is scaled if $A$ is a upper triangular matrix.
In the usual representation, a scaled linear system has the form:

$$
\left\{\begin{array}{ccc}
a_{1 r_{1}} x_{1}+ & \ldots & +a_{1 n} x_{n} \subseteq B_{1} \\
a_{2 r_{2}} x_{2}+ & \ldots & +a_{1 n} x_{n} \subseteq B_{2} \\
\vdots & & \\
a_{n r_{k}} x_{k}+ & \ldots & +a_{n n} x_{n} \subseteq B_{k}
\end{array}\right.
$$

with $r_{j} \geq 1$, and $a_{j r_{j}} \neq 0, j=1, \ldots, k$ e $r_{1}<r_{2}<\ldots<r_{k}$. For a scaled system we have three situations:

I - The last equation is of type

$$
0 x_{1}+\ldots+0 x_{n} \subseteq B_{p} \text { with } 0 \notin B_{p}
$$

In this case $S$ is impossible.
II - There is no equation of type (I) and $p=n$.
III - There is no equation of type (I) and $p<n$.
Suppose $S$ of type (II). Then we have a situation

$$
\left\{\begin{array}{lll}
a_{11} x_{1}+ & \ldots & +a_{1 n} x_{n} \subseteq B_{1} \\
a_{23} x_{3}+ & \ldots & +a_{2 n} x_{n} \subseteq B_{2} \\
\vdots & & \\
& \ldots & +a_{n n} x_{n} \subseteq B_{n}
\end{array}\right.
$$

with $a_{i i} \neq 0$ for all $i=1, \ldots, n$. Getting $x_{1}, \ldots, x_{n}$ recursively by suitable choices

$$
\begin{cases}x_{n} & \in a_{n n}^{-1} B_{n} \\ x_{k} & \in a_{k k}^{-1} B_{k}-\left[a_{k k}^{-1} a_{k(k+1)}\right] x_{k+1}-\ldots-\left[a_{k k}^{-1} a_{k n}\right] x_{n} \text { for } k=n-1, \ldots, 1\end{cases}
$$

we have a weak solution of the system $S$ (this solution is weak basically because $a_{k k}^{-1} a_{k n}$ is not a singleton in general). The same reasoning shows that for a scaled system of type (III), we can find a parametric weak solution for the system.

Example 3.3.6. Consider $n=2$ and the system over an associative superfield $F$,

$$
\left(\begin{array}{ll}
a & b \\
0 & c
\end{array}\right)\binom{x}{y} \subseteq\binom{D_{1}}{D_{2}}
$$

or in our "coloquial" representation, the system

$$
\left\{\begin{aligned}
a x+b y & \subseteq D_{1} \\
c y & \subseteq D_{2}
\end{aligned} \quad \text { with } a, c \neq 0 .\right.
$$

Since $D_{2} \subseteq c\left(c^{-1} D_{2}\right)$, for all $d_{2} \in D_{2}$ there exist $z \in c^{-1} D_{2}$ with $d_{2} \in c z$. Pick $y_{0}=z$. So we have

$$
c y_{0} \subseteq c\left(c^{-1} D_{2}\right) \text { with } c y_{0} \cap D_{2} \neq \emptyset .
$$

Hence we get $y_{0} \in c^{-1} D_{2}$ and we can choose $x_{0} \in a^{-1} D_{1}-b y_{0}$ in order to obtain

$$
a x_{0}+b y_{0} \subseteq a\left(a^{-1} D_{1}-b y_{0}\right)+b y_{0} \subseteq a a^{-1} D_{1}-a a^{-1} b y_{0}+b y_{0} \text { and } a x_{0}+b y_{0} \cap D_{1} \neq \emptyset .
$$

Of course, linear systems over superfields yields to more flexibility than linear systems over fields. It is "easier" to get a weak solution of a linear systems over superfields than over a field as we see in the Example below.

Example 3.3.7. Let $F=\mathbb{Q} /{ }_{m} \mathbb{Q}^{* 2}$. We have $2,5 \notin D(\langle 1,1\rangle)$, because $2=1 \cdot 1^{2}+1 \cdot 1^{2}$ and $5=1 \cdot 2^{2}+1 \cdot 1^{2}$. Then $\overline{2}, \overline{5} \in \overline{1}+\overline{1}$ and the system

$$
\left\{\begin{aligned}
x+y & \subseteq\{\overline{1}\} \\
y & \subseteq\{\overline{5}\}
\end{aligned}\right.
$$

over $F$ has at least a weak solution $x=y=1$.
Definition 3.3.8 (Elementary Operations). Let $A \in M_{m \times n}(F)$. The elementary operations are:
$I$ - Permute lines i e j; which will be indicated by $L_{i} \leftrightarrow L_{j}$;
II - Multiply each coefficient of a line $i$ by an element $\lambda \neq 0$ in $F$; which will be indicated by $L_{i} \leftarrow \lambda L_{i} ;$

III - Sum line $i$ with line $j$ and keep the result in line $i$; which will be indicated by $L_{i} \leftarrow L_{i}+L_{j}$.
Of course, given a linear system $A x \subseteq B$, we generate more than one system after the application of a sequence of elementary operations on the matrices $A$ and $B$. We denote the systems obtained by a set of systems $A x \subseteq B$ (with $A \subseteq M_{m \times n}(F), B \subseteq M_{m \times 1}(F)$ ) after the sequence of elementary operations $O=\left\{o_{1}, \ldots, o_{n}\right\}$ by $(A x \subseteq B)^{O}$.

The elementary operations defined above could be described in terms of matrix multiplication (as we usually do for fields). For example, considering the matrix $A \in M_{2 \times 2}(F)$ given by

$$
A=\left(\begin{array}{ll}
a & b \\
c & d
\end{array}\right)
$$

the application of $L_{1} \leftrightarrow L_{2}$ is just

$$
\left(\begin{array}{ll}
0 & 1 \\
1 & 0
\end{array}\right)\left(\begin{array}{ll}
a & b \\
c & d
\end{array}\right)=\left(\begin{array}{ll}
c & d \\
a & b
\end{array}\right)
$$

If $R$ is a proto-full superfield, to realize an elementary operation on the system $A x \subseteq B$ is
equivalent to multiply $A$ and $B$ by an elementary matrix ${ }^{4} E \in M_{m \times m}(F)$ in order to obtain the system $(E A) x \subseteq(E B)$.
element-sol
Lemma 3.3.9. Let $A x \subseteq B$ be a set of Linear systems and $O=\left\{o_{1}, \ldots, o_{n}\right\}$. Then every solution of a system in $A x \subseteq B$ is a solution in some system in $(A x \subseteq B)^{O}$. If $F$ is full, then every weak solution of a system in $A x \subseteq B$ is a weak solution in some system in $(A x \subseteq B)^{O}$.

Proof. We only need to deal with the elementary operations. Consider a system

$$
\left(\begin{array}{cccc}
a_{11} & a_{12} & \ldots & a_{1 n} \\
\vdots & \vdots & \ddots & \vdots \\
a_{m 1} & d_{m 2} & \ldots & a_{m n}
\end{array}\right)\left(\begin{array}{c}
x_{1} \\
x_{2} \\
\vdots \\
x_{n}
\end{array}\right) \subseteq\left(\begin{array}{c}
B_{1} \\
B_{2} \\
\vdots \\
B_{n}
\end{array}\right)
$$

We already know that operations (I) and (II) preserves solutions. Now consider without loss of generalization the elementary operation $L_{1} \leftarrow L_{1}+L_{2}$. Then we arrive at the set of systems

$$
\left(\begin{array}{cccc}
a_{11}+a_{21} & a_{12}+a_{22} & \ldots & a_{1 n}+a_{2 n} \\
\vdots & \vdots & \ddots & \vdots \\
a_{m 1} & d_{m 2} & \ldots & a_{m n}
\end{array}\right)\left(\begin{array}{c}
x_{1} \\
x_{2} \\
\vdots \\
x_{n}
\end{array}\right) \subseteq\left(\begin{array}{c}
B_{1}+B_{2} \\
B_{2} \\
\vdots \\
B_{n}
\end{array}\right)
$$

If we get $d_{1}, \ldots, d_{n}$ with

$$
\left(\begin{array}{cccc}
a_{11} & a_{12} & \ldots & a_{1 n} \\
\vdots & \vdots & \ddots & \vdots \\
a_{m 1} & d_{m 2} & \ldots & a_{m n}
\end{array}\right)\left(\begin{array}{c}
d_{1} \\
d_{2} \\
\vdots \\
d_{n}
\end{array}\right) \subseteq\left(\begin{array}{c}
B_{1} \\
B_{2} \\
\vdots \\
B_{n}
\end{array}\right)
$$

in particular

$$
\left(a_{11} d_{1}+\ldots+d_{1 n} d_{n}\right)+\left(a_{21} d_{1}+\ldots+d_{2 n} d_{n}\right) \subseteq B_{1}+B_{2}
$$

and

$$
\left(\begin{array}{cccc}
a_{11}+a_{21} & a_{12}+a_{22} & \ldots & a_{1 n}+a_{2 n} \\
\vdots & \vdots & \ddots & \vdots \\
a_{m 1} & d_{m 2} & \ldots & a_{m n}
\end{array}\right)\left(\begin{array}{c}
d_{1} \\
d_{2} \\
\vdots \\
d_{n}
\end{array}\right) \subseteq\left(\begin{array}{c}
B_{1}+B_{2} \\
B_{2} \\
\vdots \\
B_{n}
\end{array}\right)
$$

proving (after induction) that every solution of $A x \subseteq B$ is a solution $(A x \subseteq B)^{O}$. For the weak solution part, suppose $F$ is full and $d_{1}, \ldots, d_{n}$ with

$$
\left(\begin{array}{cccc}
a_{11} & a_{12} & \ldots & a_{1 n} \\
\vdots & \vdots & \ddots & \vdots \\
a_{m 1} & d_{m 2} & \ldots & a_{m n}
\end{array}\right)\left(\begin{array}{c}
d_{1} \\
d_{2} \\
\vdots \\
d_{n}
\end{array}\right) \cap\left(\begin{array}{c}
B_{1} \\
B_{2} \\
\vdots \\
B_{n}
\end{array}\right) \neq \emptyset
$$

[^9]In particular

$$
\left[\left(a_{11}+a_{21}\right) d_{1}+\ldots+\left(a_{1 n}+a_{2 n}\right) d_{n}\right]=\left(a_{11} d_{1}+\ldots+d_{1 n} d_{n}\right)+\left(a_{21} d_{1}+\ldots+d_{2 n} d_{n}\right) \cap\left(B_{1}+B_{2}\right) \neq \emptyset,
$$

and

$$
\left(\begin{array}{cccc}
a_{11}+a_{21} & a_{12}+a_{22} & \ldots & a_{1 n}+a_{2 n} \\
\vdots & \vdots & \ddots & \vdots \\
a_{m 1} & d_{m 2} & \ldots & a_{m n}
\end{array}\right)\left(\begin{array}{c}
d_{1} \\
d_{2} \\
\vdots \\
d_{n}
\end{array}\right) \cap\left(\begin{array}{c}
B_{1}+B_{2} \\
B_{2} \\
\vdots \\
B_{n}
\end{array}\right) \neq \emptyset .
$$

proving (after induction) that every weak solution of $A x \subseteq B$ is weak a solution $(A x \subseteq B)^{O}$.
Given a system $A x \subseteq B$, we can obtain a set of scaled systems $(A x \subseteq B)^{\text {scaled }}$, after a finite sequence of elementary operations in the same way as usual. Unfortunately, despite the result obtained in Lemma 3.3 .9 we do not know if solutions of $(A x \subseteq B)^{\text {scaled }}$ are solutions of $A x \subseteq B$.

From now on, given a system $A x \subseteq B$, to solve $A x \subseteq B$ will have the meaning to find a weak solution of $A x \subseteq B$, and a $n \times n$ system will mean a system $A x \subseteq B$ with $A \in M_{n \times n}(F)$ (and $B \in M_{n \times 1}(F)$ ).

Definition 3.3.10. Let $A=\left(a_{i j}\right) \in M_{n}(F)$ and denote $A_{i}=\left(a_{i 1}, \ldots, a_{i n}\right)$ the $i$-th row and $A^{j}=$ $\left(a_{1 j}, \ldots, a_{n j}\right)$ the $j$-th column We say that $A_{i}$ is a linear combination of $\left\{A_{1}, \ldots, A_{i-1}, A_{i+1}, \ldots, A_{n}\right\}$ if there exist $a_{1}, \ldots, a_{i-1}, a_{i+1}, \ldots, a_{n} \in F$ such that
$A_{i} \cap\left[\left(\sum_{j=1}^{r_{1}} \lambda_{j 1}\right) A_{1}+\ldots+\left(\sum_{j=1}^{r_{i-1}} \lambda_{j(i-1)}\right) A_{i-1}+\left(\sum_{j=1}^{r_{i+1}} \lambda_{j(i+1)}\right) A_{i+1}+\ldots+\left(\sum_{j=1}^{r_{n}} \lambda_{j n}\right) A_{n}\right] \neq \emptyset$. scal4
Lemma 3.3.11. Let $F$ be a superfield and $A \in M_{n}(F)$ a upper triangular matrix. Then $A$ is invertible iff $0 \notin \operatorname{det}(A)$.

Proof. Let

$$
A=\left(\begin{array}{cccc}
a_{11} & a_{12} & \ldots & a_{1 n} \\
0 & a_{22} & \ldots & a_{2 n} \\
\vdots & \vdots & \ddots & \vdots \\
0 & 0 & \ldots & a_{n n}
\end{array}\right)
$$

We already know that $\operatorname{det}(A)=a_{11} a_{22} \ldots a_{n n}$. Then, we need to prove that $A$ is invertible iff $a_{i i} \neq 0$ for all $i=1,2, \ldots, n$.

Now, let $B \in M_{n}(F)$ be another upper triangular matrix, saying

$$
B=\left(\begin{array}{cccc}
b_{11} & b_{12} & \ldots & b_{1 n} \\
0 & b_{22} & \ldots & b_{2 n} \\
\vdots & \vdots & \ddots & \vdots \\
0 & 0 & \ldots & b_{n n}
\end{array}\right)
$$

We also know that

$$
A B=\left(\begin{array}{ccccc}
a_{11} b_{11} & a_{11} b_{12}+a_{12} b_{22} & a_{11} b_{13}+a_{12} b_{23}+a_{13} b_{33} & \ldots & \sum_{k=1}^{n} a_{1 k} b_{k n} \\
0 & a_{22} b_{22} & a_{22} b_{23}+a_{23} b_{33} & \ldots & \sum_{k=2}^{n} a_{2 k} b_{k n} \\
\vdots & \vdots & \vdots & \ddots & \vdots \\
0 & 0 & 0 & \ldots & a_{n n} b_{n n}
\end{array}\right)
$$

and

$$
B A=\left(\begin{array}{ccccc}
b_{11} a_{11} & b_{11} a_{12}+b_{12} a_{22} & b_{11} a_{13}+b_{12} a_{23}+b_{13} a_{33} & \cdots & \sum_{k=1}^{n} b_{1 k} a_{k n} \\
0 & b_{22} a_{22} & b_{22} a_{23}+b_{23} a_{33} & \cdots & \sum_{k=2}^{n} b_{2 k} a_{k n} \\
\vdots & \vdots & \vdots & \ddots & \vdots \\
0 & 0 & 0 & \cdots & b_{n n} a_{n n}
\end{array}\right)
$$

Then, if $a_{i i}=0$ for some $i \in\{1,2, \ldots, n\}$, we have $I_{n} \notin A B \cap B A$ for all $B \in M_{n}(F)$. This is enough to prove that $A$ cannot be invertible.

Now, suppose $a_{i i} \neq 0$ for all $i=1,2, \ldots, n$. We will choose the elements $b_{i j}$ in order to get

$$
I_{n} \in A B \cap B A
$$

First, choose $b_{i i}=a_{i i}^{-1}$. Then, considering $A B=\left(P_{i j}\right)$ and $B A=\left(Q_{i j}\right)$, we want to get $0 \in P_{i j}$ and $0 \in Q_{i j}$ for all $i \neq j$. We need to choose $b_{(n-1) n}$ in order to get

$$
\begin{aligned}
& 0 \in a_{(n-1)(n-1)} b_{(n-1) n}+a_{(n-1) n} b_{n n} \text { and } \\
& 0 \in b_{(n-1)(n-1)} a_{(n-1) n}+b_{(n-1) n} a_{n n} .
\end{aligned}
$$

Then (remember that $b_{i i}=a_{i i}^{-1}$ ) we need

$$
b_{(n-1) n} \in-\left[a_{(n-1) n} a_{(n-1)(n-1)}^{-1}\right] a_{n n}^{-1} .
$$

Then we choose $b_{n n}, b_{(n-1)(n-1)}$ and $b_{(n-1) n}$ (i,e, we complete the process for the $n$-th and $(n-1)$-th rows of $B)$.

Now, we need to choose $b_{(n-2)(n-1)}$ and $b_{(n-2) n}$ in order to get

$$
\begin{aligned}
& 0 \in a_{(n-2)(n-2)} b_{(n-2)(n-1)}+a_{(n-2)(n-1)} b_{(n-1)(n-1)}+a_{(n-2) n} b_{n(n-1)} \text { and } \\
& 0 \in a_{(n-2)(n-2)} b_{(n-2) n}+a_{(n-2)(n-1)} b_{(n-1) n}+a_{(n-2) n} b_{n n}
\end{aligned}
$$

and

$$
\begin{aligned}
& 0 \in b_{(n-2)(n-2)} a_{(n-2)(n-1)}+b_{(n-2)(n-1)} a_{(n-1)(n-1)}+b_{(n-2) n} a_{n(n-1)} \text { and } \\
& 0 \in b_{(n-2)(n-2)} a_{(n-2) n}+b_{(n-2)(n-1)} a_{(n-1) n}+b_{(n-2) n} a_{n n}
\end{aligned}
$$

Picking $b_{(n-2)(n-1)}$ and $b_{(n-2) n}$ such that

$$
\begin{gathered}
b_{(n-2) n} \in-\left[a_{(n-2)(n-2)}^{-1} a_{(n-2)(n-1)}\right] b_{(n-1) n}-\left[a_{(n-2)(n-2)}^{-1} a_{(n-2) n}\right] b_{n n} \text { and } \\
b_{(n-2)(n-1)} \in-\left[a_{(n-2)(n-2)}^{-1} a_{(n-2)(n-1)}\right] b_{(n-1)(n-1)}-\left[a_{(n-2)(n-2)}^{-1} a_{(n-2) n}\right] b_{n(n-1)}
\end{gathered}
$$

we complete the process for the $n$-th, $(n-1)$-th and $(n-2)$-th rows of $B$. Repeating this process more $n-3$ times we arrive at a matrix $B$ such that $I_{n} \in A B \cap B A$, as desired.

### 3.4 Multipolynomials

Even if the rings-like multi-algebraic structure have been studied for more than 70 years, the developments of notions of polynomials in the ring-like multialgebraic structure seems to have a more significant development only from the last decade: for instance in [41] some notion of multi polynomials is introduced to obtain some applications to algebraic and tropical geometry, in [6] a
more detailed account of variants of concept of multipolynomials over hyperrings is applied to get a form of Hilbert's Basissatz.

Here we will stay close to the perspective in [6]: let $(R,+,-, \cdot, 0,1)$ be a superring and set

$$
R[X]:=\left\{\left(a_{n}\right)_{n \in \omega} \in R^{\omega}: \exists t \forall n\left(n \geq t \rightarrow a_{n}=0\right)\right\} .
$$

Of course, we define the degree of $\left(a_{n}\right)_{n \in \omega} \neq \mathbf{0}$ to be the smallest $t$ such that $a_{n}=0$ for all $n>t$.
Now define the binary multioperations $+, \cdot: R[X] \times R[X] \rightarrow \mathcal{P}^{*}(R[X])$, a unary operation $-: R[X] \rightarrow R[X]$ and elements $0,1 \in R[X]$ by

$$
\begin{aligned}
&\left(c_{n}\right)_{n \in \omega} \in\left(a_{n}\right)_{n \in \omega}+\left(b_{n}\right)_{n \in \omega} \text { iff } \forall n\left(c_{n} \in a_{n}+b_{n}\right) \\
&\left(c_{n}\right)_{n \in \omega} \in\left(\left(a_{n}\right)_{n \in \omega} \cdot\left(b_{n}\right)_{n \in \omega} \text { iff } \forall n\left(c_{n} \in a_{0} \cdot b_{n}+a_{1} \cdot b_{n-1}+\ldots+a_{n} \cdot b_{0}\right)\right. \\
&-\left(a_{n}\right)_{n \in \omega}=\left(-a_{n}\right)_{n \in \omega} \\
& 0:=(0)_{n \in \omega} \\
& 1:=(1,0, \ldots, 0, \ldots)
\end{aligned}
$$

For convenience, we denote elements of $R[X]$ by $\alpha=\left(a_{n}\right)_{n \in \omega}$. Beside this, we denote

$$
\begin{aligned}
1 & :=(1,0,0, \ldots), \\
X & :=(0,1,0, \ldots), \\
X^{2} & :=(0,0,1,0, \ldots)
\end{aligned}
$$

etc. In this sense, our "monomial" $a_{i} X^{i}$ is denoted by $\left(0, \ldots 0, a_{i}, 0, \ldots\right)$, where $a_{i}$ is in the $i$-th position; in particular, we will denote $\underline{b}=(b, 0,0, \ldots)$ and we frequently identify $b \in R$ ams $\underline{b} \in R[X]$.

The properties stated in the Lemma below immediately follows from the definitions involving $R[X]$ :

Lemma 3.4.1. Let $R$ be a superring and $R[X]$ as above and $n, m \in \mathbb{N}$.
$a-\left\{X^{n+m}\right\}=X^{n} \cdot X^{m}$.
$b-$ For all $a \in R,\left\{a X^{n}\right\}=\underline{a} \cdot X^{n}$.
$c$ - Given $\alpha=\left(a_{0}, a_{1}, \ldots, a_{n}, 0,0, \ldots\right) \in R[X]$, with with $\operatorname{deg} \alpha \leq n$ and $m \geq 1$, we have

$$
\alpha X^{m}=\left(0,0, \ldots, 0, a_{0}, a_{1}, \ldots, a_{n}, 0,0, \ldots\right)=a_{0} X^{m}+a_{1} X^{m+1}+\ldots+a_{n} X^{m+n}
$$

$d$ - For $\alpha=\left(a_{n}\right)_{n \in \omega} \in R[X]$, with $\operatorname{deg} \alpha=t$,

$$
\{\alpha\}=a_{0} \cdot 1+a_{1} \cdot X+\ldots+a_{t} \cdot X^{t}=a_{0}+X\left(a_{1}+a_{2} X+\ldots+a_{n} X^{t-1}\right) .
$$

$e-R[X]$ is a non-associative superring. If $R$ is associative and full then $R[X]$ is an associative superring.
$f-R[X]$ is a superdomain iff $R$ is a superdomain.
$g$ - The map $a \in R \mapsto \underline{a}=(a, 0, \cdots, 0, \cdots)$ defines a full embedding $R \mapsto R[X]$.
$h$ - For an ordinary ring $R$ (identified with a strict superring), the superring $R[X]$ is naturally isomorphic to (the superring associated to) the ordinary ring of polynomials in one variable over $R$.

Lemma 3.4.1 allow us to deal with the superring $R[X]$ as usual. In other words, we can assume that for $\alpha \in R[x]$, there exists $a_{0}, a_{1}, \ldots, a_{n} \in R$ such that $\alpha=a_{0}+a_{1} X+\ldots+a_{n} X^{n}$, and then, we can work simply denoting $\alpha=f(X)$, as usual. For example, combining the definitions and all facts above we get

$$
(x-a)(x-b)=x^{2}+(a-b) x+a b=\left\{x^{2}+d x+e: d \in a-b \text { and } e \in a b\right\} .
$$

Here we have a situation similar to the matrix case: the general structure $R[X]$ will be associative only if $R$ is full.
teo-rxfull
Theorem 3.4.2. Let $R$ be an associative proto-full superring. Then $R[X]$ is a proto-full superring. Moreover, if $R$ is full then $R[X]$ is an associative superring.

Proof. In fact, we already know that $R[X]$ is a non-associative superring. To prove the desired affirmations, we deal with elements in $R[X]$ as sequences: we denote $a=\left(a_{0}, a_{1}, \ldots, a_{n}, \ldots\right) \in R[X]$, and for $n \geq 0,[a]_{n}:=a_{n}$. We extend this notation for the operations + and $\cdot$ over $R[X]$ :

$$
\begin{aligned}
{[a+b]_{n} } & :=a_{n}+b_{n} \\
{[a b]_{n} } & :=\sum_{i=0}^{n} a_{i} b_{n-i}
\end{aligned}
$$

With these notations, for all $a, b, c \in R[X]$ and all $n \geq 0$ we get

$$
\begin{aligned}
& {[(a b) c]_{n}=\sum_{i=0}^{n}\left[[a b]_{i}\right] c_{n}=\sum_{i=0}^{n}\left[\sum_{p=0}^{i} a_{p} b_{i-p}\right] c_{n}} \\
& {[a(b c)]_{n}=\sum_{i=0}^{n} a_{i}\left[[b c]_{n-i}\right]=\sum_{i=0}^{n} a_{i}\left[\sum_{p=0}^{n-i} b_{p} c_{n-i-p}\right]}
\end{aligned}
$$

If $R$ is full (proto-full), then (after some reindexation) we get $[(a b) c]_{n}=[a(b c)]_{n}\left(\left[[(a b) c]_{n}\right] \cap\right.$ $\left[[a(b c)]_{n}\right] \neq \emptyset$ ) for all $n \geq 0$, and then, $(a b) c=a(b c)$.

Remark 3.4.3. If $R$ is a full superdomain, does not hold in general that $R[X]$ is also a full superdomain. In fact, even if $R$ is a hyperfield, there are examples, e.g. $R=K, Q_{2}$, such that $R[X]$ is not a full superdomain (see [6]).

Definition 3.4.4. The structure $R[X]$ will be called polynomials in one variable over $R$. The elements of $R[X]$ will be called polynomials. We denote $R\left[X_{1}, \ldots, X_{n}\right]:=\left(R\left[X_{1}, \ldots, X_{n-1}\right]\right)\left[X_{n}\right]$.
degreelemma
Lemma 3.4.5 (Adapted from Theorem 5 of [6]). Let $R$ be a superring and $f, g \in R[X] \backslash\{0\}$.
$i$ - If $t(X) \in f(X)+g(X)$ and $f \neq-g$ then

$$
\min \{\operatorname{deg}(f), \operatorname{deg}(g)\} \leq \operatorname{deg}(t) \leq \max \{\operatorname{deg}(f), \operatorname{deg}(g)\}
$$

ii - If $R$ is a superdomain and $t(X) \in f(X) g(X)$, then $\operatorname{deg}(t)=\operatorname{deg}(f)+\operatorname{deg}(g)$. In particular, if $f_{1}(X), f_{2}(X), \ldots, f_{n}(X) \neq 0$ and $t(X) \in f_{1}(X) f_{2}(X) \ldots f_{n}(X)$, then

$$
\operatorname{deg}(t)=\operatorname{deg}\left(f_{1}\right)+\operatorname{deg}\left(f_{2}\right)+\ldots+\operatorname{deg}\left(f_{n}\right) .
$$

iii - (Partial Factorization) Let $R$ be a superdomain, $\operatorname{deg}(f)=n$ and $f \in\left(X-a_{1}\right)\left(X-a_{n}\right) \ldots(X-$ $\left.a_{p}\right)$. Then $p=n$.

Let $f(X)=a_{0}+\ldots+a_{n} X^{n}$ and $g(X)=b_{0}+\ldots+b_{m} X^{m}$ with $a_{n}, b_{m} \neq 0$. We establish the following notation: for $k \in \mathbb{N}$ with $k \leq \operatorname{deg}(f)$ we define $(f)_{k}:=a_{k}$ (the $k$-th coefficient of $f$ ).

Proof of Lemma 3.4.5. For item (i), we have

$$
f(X)+g(X)=\left(a_{0}+b_{0}\right) X+\ldots+\left(a_{n}+b_{m}\right) X^{m} .
$$

Since $f(X) \neq-g(X), 0 \notin a_{n}+b_{n}$, establishing item (i).
Now, suppose without loss of generality that $m \geq n$ and in this case, write

$$
f(X)=a_{0}+\ldots+a_{m} X^{m}
$$

with $a_{k}=0$ for $n<k \leq m$. We have $(f g)_{m+n} \in a_{n} b_{m}$ and since $R$ is a superdomain, $(f g)_{m+n} \neq 0$. This and induction proves item (ii).

For item (iii), let $g \in\left(X-a_{1}\right)\left(X-a_{n}\right) \ldots\left(X-a_{p}\right)$. By item (ii) and induction, $\operatorname{deg}(g)=p$. Then $n=\operatorname{deg}(f)=p$.

Despite the fact that $R[X]$ is not full in general, we have a powerful Lemma to get around this situation.

Lemma 3.4.6. Let $R$ be a superring and $f \in R[X]$ with $f(X)=a_{n} X^{n}+\ldots+a_{1} X+a_{0}$. Then:
$i$ - For all $b, c \in R,(b+c X) f(X)=b f(X)+c X f(X)$.
ii - For all $b, c \in R$ and all $p, q \in \omega$ with $p<q$,

$$
\left(b X^{p}+c X^{q}\right) f(X)=b X^{p} f(X)+c X^{p} f(X) .
$$

iii - For all $b, c, d \in R$ and all $p, q, r \in \omega$ with $p<q<r$,

$$
\left(b X^{p}+c X^{q}+d X^{r}\right) f(X)=b X^{p} f(X)+c X^{p} f(X)+d X^{r} f(X) .
$$

$i v$ - For all $b_{0}, \ldots, b_{m} \in R$,

$$
\begin{array}{r}
\left(b_{0}+b_{1} X+b_{2} X^{2}+\ldots+b_{m} X^{m}\right) f(X)= \\
b_{0} f(X)+\left(b_{1} X+b_{2} X^{2}+\ldots+b_{m} X^{m}\right) f(X) .
\end{array}
$$

$v$ - For all $b_{0}, \ldots, b_{m} \in R$,

$$
\begin{array}{r}
\left(b_{0}+b_{1} X+b_{2} X^{2}+\ldots+b_{m-1} X^{m-1}+b_{m} X^{m}\right) f(X)= \\
\left(b_{0}+b_{1} X+b_{2} X^{2}+\ldots+b_{m-1} X^{m-1}\right) f(X)+b_{m} X^{m} f(X) .
\end{array}
$$

vi - For all $b_{0}, \ldots, b_{m} \in R$,

$$
\begin{array}{r}
\left(b_{0}+b_{1} X+\ldots+b_{j} X^{j}+b_{j+1} X^{j+1}+\ldots+b_{m} X^{m}\right) f(X)= \\
\left(b_{0}+b_{1} X+\ldots+b_{j} X^{j}\right) f(X)+\left(b_{j+1} X^{j+1}+\ldots+b_{m} X^{m}\right) f(X) .
\end{array}
$$

In particular, if $d \in R, g(X) \in R[X]$ and $r>\operatorname{deg}(g(X))$, then

$$
\left(g(X)+d X^{r}\right) f(X)=g(X) f(X)+d X^{r} f(X) .
$$

## Proof.

i - We can suppose without loss of generality that $b, c \neq 0$. Here is convenient keep in mind that an element in $R[X]$ is a sequence of elements in $R$. Denote $b+c X=\left(b_{n}\right)_{n \in \omega} \in R[X]$ with $b_{0}=b, b_{1}=c$ and $b_{n}=0$ for all $n \geq 2$. By definition, for an element $h(X) \in R[X]$, say

$$
h(X)=e_{0}+e_{1} X+\ldots+e_{n+1} X^{n+1}=\left(e_{n}\right)_{n \in \omega} \in R[X],
$$

we have

$$
h(X) \in(b+c X) f(X) \text { iff } e_{p} \in \sum_{j=0}^{p} a_{j} b_{p-j}, p \in \omega .
$$

Since $a_{j}=0$ for all $j>n$ and $b_{j}=0$ for all $j \geq 2$, we have $e_{p}=0$ for all $p>n+1$. Moreover, by the same reason we have that $e_{0} \in a_{0} b_{0}, e_{n+1} \in a_{n} b_{1}$ and for $0<p<n+1$, that

$$
e_{p} \in \sum_{j=0}^{p} a_{j} b_{p-j}=a_{p} b_{0}+a_{p-1} b_{1} .
$$

Summarizing, we conclude that

$$
\begin{equation*}
h(X) \in(b+c X) f(X) \text { iff } e_{0} \in a_{0} b_{0}, e_{n+1} \in a_{n} b_{1} \text { and } e_{p} \in a_{p} b_{0}+a_{p-1} b_{1} \text { for } 0<p<n+1 . \tag{}
\end{equation*}
$$

On the other hand, we have that

$$
\begin{align*}
b f(X)+c X f(X) & =b\left[a_{n} X^{n}+\ldots+a_{1} X+a_{0}\right]+c X\left[a_{n} X^{n}+\ldots+a_{1} X+a_{0}\right] \\
& =\left(a_{n} b X^{n}+\ldots+a_{1} b X+a_{0} b\right)+\left(a_{n} c X^{n+1}+\ldots+a_{1} c X^{2}+a_{0} X\right) \\
& =a_{n} c X^{n+1}+\left(a_{n} b+a_{n-1} c\right) X^{n}+\ldots+\left(a_{2} b+a_{1} c\right) X^{2}+\left(a_{1} b+a_{0} c\right) X+a_{0} b . \tag{**}
\end{align*}
$$

Joining (*) and (**) we conclude that

$$
h(X) \in(b+c X) f(X) \text { iff } h(X) \in b f(X)+c X f(X) .
$$

ii - Just use the same reasoning of item (i).
iii - Using distributivity, item (i) and (ii) we conclude that

$$
\begin{equation*}
\left(b X^{p}+c X^{q}+d X^{r}\right) f(X) \subseteq b X^{p} f(X)+c X^{p} f(X)+d X^{r} f(X)=\left(b X^{p}+c X^{p}\right) f(X)+d X^{r} f(X * *) . \tag{***}
\end{equation*}
$$

Now with $(* * *)$ on hand, just proceed with the same reasoning of $(*)$ and $(* *)$ to obtain the desired.
iv - This is an immediate consequence of item (iii) and a convenient induction.
v - This is an immediate consequence of item (iii) and a convenient induction.
vi - This is just the combination of previous items.
euclid
Theorem 3.4.7 (Euclid's Division Algorithm (3.4 in [11)). Let $K$ be a superfield. Given polynomials $f(X), g(X) \in K[X]$ with $g(X) \neq 0$, there exists $q(X), r(X) \in K[X]$ such that $f(X) \in$ $q(X) g(X)+r(X)$, with $\operatorname{deg} r(X)<\operatorname{deg} g(X)$ or $r(X)=0$.

Proof. This is a generalized version of Theorem 3.4 in [11, which states Euclid's Algorithm for hyperfields. Write

$$
\begin{aligned}
& f(X)=a_{n} X^{n}+\cdots+a_{1} X+a_{0} \\
& g(X)=b_{m} X^{m}+\cdots+b_{1} X+b_{0}
\end{aligned}
$$

with $a_{n}, b_{m} \neq 0$ and let $b_{m}^{-1} \in K$ be an element satisfying $1 \in b_{m} \cdot b_{m}^{-1}$.
We proceed by induction on $n$. Note that if $m \geq n$, then is sufficient take $q(X)=0$ and $r(X)=f(X)$, so we can suppose $m \leq n$. If $m=n=0$, then $f(X)=a_{0}$ and $g(X)=b_{0}$ are both non zero constants, so is sufficient take $q(X) \in a_{0} \cdot b_{0}^{-1}$ and $r(X)=0$.

Now, suppose $n \geq 1$. Then, since $0 \in a-a$, there exist some $t(X) \in f(X)-a_{n} b_{m}^{-1} X^{n-m} g(X)$ with $\operatorname{deg} t(X)<n$. So, by induction hypothesis,

$$
t(X) \in q(X) g(X)+r(X) \text { for some } q(X), r(X) \in R[X] \text { with } \operatorname{deg} r(X)<\operatorname{deg} g(X) \text { or } r(X)=0 .
$$

Therefore, $\operatorname{deg} t(X)=\operatorname{deg} q(X)+m$ and since $f(X) \in t(X)+a_{n} b_{m}^{-1} X^{n-m} g(X)$, we have

$$
\begin{aligned}
f(X) & \in t(X)+a_{n} b_{m}^{-1} X^{n-m} g(X) \\
& \subseteq q(X) g(X)+a_{n} b_{m}^{-1} X^{n-m} g(X)+r(X) .
\end{aligned}
$$

But since $\operatorname{deg} q(X)=\operatorname{deg} t(X)-m<n-m$, we have (see Lemma 3.4.6 (vi)) that

$$
\left[q(X)+a_{n} b_{m}^{-1} X^{n-m}\right] g(X)=q(X) g(X)+a_{n} b_{m}^{-1} X^{n-m} g(X) .
$$

So there exist some $q^{\prime}(X) \in q(X)+a_{n} b_{m}^{-1} X^{n-m}$ with $f(X) \in q^{\prime}(X) g(X)+r(X)$ and $\operatorname{deg} r(X)<$ $\operatorname{deg} g(X)$ or $r(X)=0$, completing the proof.

## Remark 3.4.8.

$i$ - Note that the polynomials $q$ and $r$ of Theorem 3.4.7 are not unique in general: if $f \in g q+r$, then $f \in g(q+1-1)+r$ and $f \in g q+(r+1-1)$, then, if $\{0\} \neq 1-1$, we have many $q$ 's and $r$ 's.
ii - However, if $R$ is a ring, then Theorem 3.4.7 provide the usual Euclid Algorithm, with the uniqueness of the quotient and remainder.
teoPID
Theorem 3.4.9 (Adapted from Theorem 6 of [6]). Let $F$ be a full associative superfield. Then $F[X]$ is a principal ideal superdomain.

Proof. Let $I$ be a ideal of $F[X]$. If $I=0$ then $I=\langle 0\rangle$ and if there is some $a \in F \backslash\{0\}$ with $a \in I$, then $I=F[X]=\langle 1\rangle$ (because $F$ is full).

Now let $p(X) \in I$ be a polynomial with minimal degree $m \geq 1$. Let $f(X) \in I$ be another polynomial. By Euclid's Algorithm, there exists $q(X), r(X) \in F[X]$ with $f(X) \in p(X) q(X)+r(X)$ and $r(X)=0$ or $\operatorname{deg}(r)<\operatorname{deg}(p)=m$. Since $f, p \in I$ and $r(X) \in f(X)-p(X) q(X)$, we have
$r \in I$. Note that by the minimality of $m$, all nonzero polynomial in $f(X)-p(X) q(X)$ has degree at least $m$. If $r \neq 0$ then

$$
\min \{\operatorname{deg} f, \operatorname{deg}(p)+\operatorname{deg}(q)\} \leq \operatorname{deg} r \leq \max \{\operatorname{deg} f, \operatorname{deg}(p)+\operatorname{deg}(q)\} .
$$

In particular $\operatorname{deg}(r) \geq m$ (because $\operatorname{deg}(f) \leq m$ ), contradicting $\operatorname{deg}(r)<m$. Hence $r=0$ and $I=\langle p\rangle$. In particular, $I=F[X] \cdot p(X)$.

### 3.5 Evaluation and Roots

Let $R, S$ be superrings and $h: R \rightarrow S$ be a morphism. Then $h$ extends naturally to a morphism in the proto-superrings multipolynomials $h^{X}: R[X] \rightarrow S[X]$ :

$$
\left(a_{n}\right)_{n \in \mathbb{N}} \in R[X] \mapsto\left(h\left(a_{n}\right)\right)_{n \in \mathbb{N}} \in S[X]
$$

Now let $s \in S$. We define the $h$-evaluation of $s$ at $f(X) \in R[X]$ with $f(X)=a_{0}+a_{1} X+\ldots+$ $a_{n} X^{n}$ by

$$
f^{h}(s)=e v^{h}(s, f):=\left\{s^{\prime} \in S: s^{\prime} \in h\left(a_{0}\right)+h\left(a_{1}\right) s+h\left(a_{2}\right)\left(s^{2}\right)+\ldots+h\left(a_{n}\right)\left(s^{n}\right)\right\} .
$$

In order to easy our presentation, we just denote $a b^{n}:=a\left(b^{n}\right)$. We define the $h$-evaluation for a subset $I \subseteq S$ by

$$
f^{h}(I)=\bigcup_{s \in I} f^{h}(s)
$$

In particular if $S \supseteq R$ are superrings and $\alpha \in S$, we have the evaluation of $\alpha$ at $f(X) \in R[X]$ by

$$
f(\alpha, S)=e v(\alpha, f, S)=\left\{b \in S: b \in a_{0}+a_{1} \alpha+a_{2} \alpha^{2}+\ldots+a_{n} \alpha^{n}\right\} \subseteq S .
$$

Note that the evaluation depends on the choice of $S$. When $S=R$ we just denote $f(\alpha, R)$ by $f(\alpha)$.

A root of $f$ in $S$ is an element $\alpha \in S$ such that $0 \in e v(\alpha, f, S)$. In this case we say that $\alpha$ is $S$-algebraic over $R$. An effective root of $f$ in $S$ is an element $\alpha \in S$ such that $f \in(X-\alpha) \cdot g(X)$ for some $g(X) \in R[X]$. A superring $R$ is algebraically closed if every non constant polynomial in $R[X]$ has a root in $R$.

Observe that, if $F$ is a field, the evaluation of $F[X]$ as a ring coincide with the usual evaluation, and, of course, root and effective roots are the same thing. Therefore, if $F$ is algebraically closed as hyperfield and superfield, then will be algebraically closed in the usual sense.

Remark 3.5.1. The expansion of the above field-theoretical concepts to the multialgebraic theory of superfields (hyperfields, in particular) brings new phenomena:
$i$ - (Polynomials can have infinite roots): Let $F$ be a infinite pre-special hyperfield (24]). Then $F$ has characteristic $0, a^{2}=1$ for all $a \neq 0$ so the polynomial $f(X)=X^{2}-1$ has infinite roots (i.e, $0 \in e v(f, \alpha)$ for all $\alpha \in \dot{F})$.
ii- (Finite hyperfields can be algebraically closed). The hyperfield $K=\{0,1\}$ is algebraically closed. In fact, if $p(X)=a_{0}+a_{1} X+a_{2} X^{2}+\ldots+a_{n} X^{n} \in K[X]$, with $a_{n} \neq 0$, then $0 \in p(0)$ (if $a_{0}=0$ ) or $p(1)=K$, since $1+1=\{0,1\}$.

We have good results concerning irreducibly (see for instance, Theorem 3.5.4 below). These results are the key to the development of superfields extensions, which leads us to some kind of algebraic closure.

Definition 3.5.2 (Irreducibility). Let $R$ be a superfield and $f, d \in R[X]$. We say that divides $f$ if and only if $f \in\langle d\rangle$, and denote $d \mid f$. We say that $f$ is irreducible if $\operatorname{deg} f \geq 1$ and $u \mid f$ for some $u \in R[X]$ (i.e, $f \in\langle u\rangle$ ), then $\langle f\rangle=\langle u\rangle$.

Theorem 3.5.3. Let $F$ be a full associative superfield and $p(X) \in F[X]$ be an irreducible polynomial. Then $\langle p(X)\rangle$ is a maximal ideal.

Proof. Let $p(X)$ be irreducible and $I \subseteq F[X]$ an ideal with $\langle p(X)\rangle \subseteq I$. By Theorem 3.4.9.

$$
I=\langle f(X)\rangle=F[X] \cdot f(X)
$$

for some $f(X) \in F[X]$. Since $p(X) \in I=\langle f(X)\rangle$, then $p(X)=f(X) g(X)$ for some $g(X) \in F[X]$. Since $p(X)$ is irreducible, either $f(X)$ or $g(X)$ is a constant polynomial. If $f(X)$ is constant, then $I=F[X]$, and if $g(X)$ is constant, $I=\langle p(X)\rangle$, which proves that $\langle p(X)\rangle$ is maximal.

If $F$ is not full, we cannot prove that $\langle p(X)\rangle$ is a maximal ideal. But we still have that $F[X] /\langle p\rangle$ is a superfield.
lemquadext
Theorem 3.5.4. Let $F$ be a superfield and $p \in F[X]$ be an irreducible polynomial. Then $F[X] /\langle p\rangle$ is a superfield.

Proof. Let $p(X)=d_{0}+a_{1} X+\ldots+a_{n+1} X^{n+1}$. Note that

$$
\begin{align*}
F(p(X)):=F[x] /\langle p\rangle & =\left\{\left[a_{0}+a_{1} X+\ldots+a_{n} X^{n}\right]: a_{0}, \ldots, a_{n} \in F\right\} \\
& =\left\{[f(X)]: f(X)=a_{0}+a_{1} X+\ldots+a_{r} X^{r} \text { with } a_{0}, \ldots, a_{r} \in F, r \leq n\right\} .
\end{align*}
$$

Let $f(X)=a_{0}+a_{1} X+\ldots+a_{r} X^{r}$ and $g(X)=b_{0}+b_{1} X+\ldots+b_{s} X^{s}$ with and suppose

$$
[0] \in[f(X)][g(X)] .
$$

There exist

$$
h(X) \in(f(X) g(X)) \cap\langle p(X)\rangle .
$$

Since $F$ is a superdomain, every nonzero polynomial in $\langle p\rangle$ has degree at least $n+1=\operatorname{deg}(p)$. Now get a nonzero element in $[t(x)] \in[f(X)][g(X)]$. Using Equation 3.2 we have $t(X) \in f(X) g(X)$ with $\operatorname{deg}(t) \leq n$. Then $h(X)=0$ and $0 \in f(X) g(X)$, which imply $f(X)=0$ or $g(X)=0$ (because $F[X]$ is a superdomain). Then $[F(X)]=0$ or $[g(X)]=[0]$, proving that $F[p(X)]$ is a superdomain (and then, $\langle p(X)\rangle$ is strongly prime).

Now we prove that $F[p(X)]$ is a superfield, i.e, that for all nonzero $[f(X)] \in F[p(X)]$, there exist a nonzero $[g(X)] \in F[p(X)]$ with $[1] \in[f(X)][g(X)]$. We proceed by induction on $n=\operatorname{deg}(f(X))$.

If $n=0$, then $f(X)=a$ for some $a \in \dot{F}$, and there exist $a^{-1} \in \dot{F} \dot{H}^{5}$ with $1 \in a \cdot a^{-1}$, and then $[1] \in[f(X)]\left[a^{-1}\right]$. If $n=1$, then $f(X)=a X+b, a, b \in F(a \neq 0)$. By Euclid's Algorithm, there exists $q(X), r(X)$ with $p(X) \in f(X) q(X)+r(X)$ with $r(X)=0$ or $\operatorname{deg}(r(X))<\operatorname{deg}(f(X))$. Since $p(X)$ is irreducible, $r(X) \neq 0$ and $r(X)=d \in \dot{F}$. Moreover for some $d^{-1} \in \dot{F}$ with $1 \in d \cdot d^{-1}$ we

[^10]have
\[

$$
\begin{aligned}
p(X) & \in f(X) q(X)+d \Rightarrow[0] \in[f(X)][q(X)]+[d] \Rightarrow-[d] \in[f(X)][q(X)] \\
& \Rightarrow\left[d d^{-1}\right] \subseteq[f(X)]\left(-\left[d^{-1}\right][q(X)]\right) \Rightarrow[1] \in[f(X)]\left(-\left[d^{-1}\right][q(X)]\right),
\end{aligned}
$$
\]

and then, there exist $[t(X)] \in-\left[d^{-1}\right][q(X)]$ with $[1] \in[f(X)][t(X)]$.
Now, suppose by induction that all polynomial of degree at most $n$ has an inverse and let $f(X) \in F[X]$ with $\operatorname{deg}(f(X))=n+1$. By Euclid's Algorithm, there exists $q(X), r(X)$ with $p(X) \in f(X) q(X)+r(X)$ with $r(X)=0$ or $\operatorname{deg}(r(X))<\operatorname{deg}(f(X))$ and since $p(X)$ is irreducible, we have $r(X) \neq 0$. By induction hypothesis, there exist $g(X) \in F[X]$ with $[1] \in[r(X)][g(X)]$. Then

$$
\begin{aligned}
p(X) \in f(X) q(X)+r(X) & \Rightarrow[0] \in[f(X)][q(X)]+[r(X)] \\
& \Rightarrow[r(X)] \in-[f(X)][q(X)] \\
& \Rightarrow[r(X)][g(X)] \subseteq-[f(X)][q(X)][g(X)] \\
& \Rightarrow 1 \in[r(X)][g(X)] \subseteq[f(X)](-[q(X)][g(X)]),
\end{aligned}
$$

then there exist $[t(X)] \in-[q(X)][g(X)]$ with $[1] \in[f(X)][t(X)]$, completing the proof.
Using Theorem 3.5.4 we obtain an algorithm to determine the invertible elements in $F[p(X)]$ particularly useful in the field case:

Corollary 3.5.5. Let $F$ be a field and $p(X) \in F[X]$ be an irreducible polynomial. If $f(X) \neq 0$ and $p(X)=f(X) q(X)+r(X)$ with $r(X) \neq 0$, then

$$
[f(X)]^{-1}=-[q(X)][r(X)]^{-1} \in F[p(X)] .
$$

Definition 3.5.6. Let $F$ be a superfield and $p(X) \in F[X]$ be an irreducible polynomial. We denote $F(p):=F(p(X))=F[X] /\langle p(X)\rangle$.

Lemma 3.5.7. Let $F$ be a superfield and $p(X) \in F[X]$ be an irreducible polynomial. Denote $\bar{X}=\lambda$ and let $f \in F(p)$ with $f=\overline{a_{n}} \lambda^{n}+\ldots+\overline{a_{1}} \lambda+\bar{a}_{0}$. Then:
$i$ - For all $b, c \in F,(\bar{b}+\bar{c} \lambda) f=\bar{b} f+\bar{c} \lambda f$.
$i i-$ For all $b_{0}, \ldots, b_{m} \in F$,

$$
\begin{gathered}
\left(\overline{b_{0}}+\overline{b_{1}} \lambda+\ldots+\overline{b_{j}} \lambda^{j}+\overline{b_{j+1}} \lambda^{j+1}+\ldots+\overline{b_{m}} \lambda^{m}\right) f= \\
\left(\overline{b_{0}}+\overline{b_{1}} \lambda+\ldots+\overline{b_{j}} \lambda^{j}\right) f+\left(\overline{b_{j+1}} \lambda^{j+1}+\ldots+\overline{b_{m}} \lambda^{m}\right) f .
\end{gathered}
$$

In particular, if $d \in F, g \in F(p)$ with $g=\overline{b_{0}}+\overline{b_{1}} \lambda+\overline{b_{2}} \lambda^{2}+\ldots+\overline{b_{m}} \lambda^{m}$ and $r>m$, then

$$
\left(g+\bar{d} \lambda^{r}\right) f=g f+\bar{d} \lambda^{r} f
$$

Proof. Similar to Lemma 3.4.6.
Theorem 3.5.8. Let $F$ be a superfield and $p(X) \in F[X]$ be a polynomial of degree greater or equal to 1. Then there exist a superfield $L$ such that $F \subseteq L, F$ is a sub superfield of $L$ (i.e, the inclusion $F \hookrightarrow L$ is a full morphism) and $p(X)$ has a root.

Proof. It is enough to show the result for $p(X)$ irreducible. In this case, the ideal $\langle p(X)\rangle \subseteq F[X]$ is maximal and $K^{\prime}=F[X] /\langle p(X)\rangle$ is a superfield. If we consider the canonical injection $\iota: F \rightarrow$ $F[X] /\langle p\rangle$ given by $a \mapsto \bar{a}$, we have a full morphism (basically because $F \hookrightarrow F[X]$ is full). Putting $F^{\prime}=\iota(F)$ we have that $F \cong F^{\prime}, F^{\prime} \hookrightarrow L$ is a full morphism and the polynomial $p^{\iota}$ (given by the application of $\iota$ in each coefficient) has a root $\bar{x}$.

Next, let $K=F \cup X$ for some $X$ of cardinality $K^{\prime} \backslash F^{\prime}$. We construct a bijection $\varphi: K \rightarrow K^{\prime}$ which restrict ton $F$ is equal to $\iota$. This bijection transport the structure of superfield for $K$ (in the obvious way), in order to get an extension $K \mid F$ such that $f$ has a root $\varphi^{-1}(\bar{x})$.

Corollary 3.5.9. Let $F$ be a superfield and $f \in F[X]$ be a polynomial with $n=\operatorname{deg}(f) \geq 1$. Then there exist a non-associative superfield $L$ such that $F \subseteq L$ and $f$ has at least $n$ roots.

Corollary 3.5.10. Let $F$ be a superfield and $f_{1}, \ldots, f_{n} \in F[X]$ be polynomials with $1 \leq \operatorname{deg}\left(f_{j}\right)=$ $r_{j}, j=1, \ldots, n$. Then there exist a superfield $L$ such that $F \subseteq L$ and each $f_{j}$ has at least $r_{j}$ roots.

### 3.6 Extensions

We have some possibilities to consider in order to define the notion of extension for superfields:
extension
Definition 3.6.1 (Extensions). Let $F$ and $K$ be superfields.
$i$ - We say that $K$ is a proto superfield extension (or just a proto extension) of $F$, notation $\left.K\right|_{p} F$, if $F \subseteq K$.
ii - We say that $K$ is a superfield extension (or just an extension) of $F$, notation $K \mid F$ if $F \subseteq K$ and the inclusion map $F \hookrightarrow K$ is a superfield morphism.
iii - We say that $K$ is a full superfield extension (or just a full extension) of $F$, notation $\left.K\right|_{f} F$ if $F \subseteq K$ and the inclusion map $F \hookrightarrow K$ is a full superfield morphism.

## Example 3.6.2.

$i$ - Of course, all full extension is an extension and all extension is a proto extension.
ii - We have $K \subseteq Q_{2}$ but the inclusion map $K \hookrightarrow Q_{2}$ is not a morphism. Then we have a proto extension $\left.Q_{2}\right|_{p} K$ that is not an extension.
iii - For $p, q$ prime integers with $q \geq p$ we have an inclusion morphism $H_{p} \hookrightarrow H_{q}$, but this morphism is not full. Then we have an extension $H_{q} \mid H_{p}$ that is not a full extension.
iv - Let $F$ be a superfield, $p \in F[X]$ an irreducible polynomial and $F(p)=F[X] /\langle p\rangle$ be the superfield built in Theorem 3.5.8. Then we have a full morphism $F \hookrightarrow F(p)$ so we have a full extension $\left.F(p)\right|_{f} F$.
$v$ - Let $F, K$ be fields such that $F \subseteq K$. Then the field extension $K \mid F$ satisfy all conditions in Definition 3.6.1.

The result below justify a deeper look at full superfield extensions.
unicityext
Theorem 3.6.3. Let $\left.K_{1}\right|_{f} F$ and $\left.K_{2}\right|_{f} F$ be full superfield extensions and suppose that $\gamma \in K_{1} \cap K_{2}$. Then

$$
F\left[\gamma, K_{1}\right]=F\left[\gamma, K_{2}\right] .
$$

Proof. Suppose first that $\left.K_{2}\right|_{f} K_{1}$ is a full extension. Then for all $f \in F[X], e v\left(f, K_{1}\right)=e v\left(f, K_{2}\right)$, so $F\left[\gamma, K_{1}\right]=F\left[\gamma, K_{2}\right]$.

Now, for the general case just note that $\left.K_{1}\right|_{f}\left(K_{1} \cap K_{2}\right)$ and $\left.K_{2}\right|_{f}\left(K_{1} \cap K_{2}\right)$. Then

$$
F\left[\gamma, K_{1}\right]=F\left[\gamma, K_{1} \cap K_{2}\right]=F\left[\gamma, K_{2}\right]
$$

Definition 3.6.4 (Algebraic Extensions). We say that a proto extension $\left.K\right|_{p} F$ is algebraic if all element $\alpha \in K$ is $K$-algebraic over $F$. We denote the same for extensions and full extensions.

Definition 3.6.5 (Linear Independency, Basis, Degree). Let $\left.K\right|_{p} F$ be a proto extension and $I \subseteq K$. We say that $I$ is $F$-linearly independent if for all distinct $\lambda_{1}, \ldots, \lambda_{n} \in I, n \in \mathbb{N}$, the following hold:

$$
\text { If } 0 \in a_{1} \lambda_{1}+\ldots+a_{n} \lambda_{n} \text { then } a_{1}=\ldots=a_{n}=0
$$

and $I$ is F-linearly dependent if it is not F-linearly independent. We say that $I$ is a F-basis of $K$ if $I$ is linearly independent and $K$ is generated by $I$, i.e,

$$
K=\bigcup_{n \geq 0}\left\{\sum_{i=0}^{n} a_{i} \lambda_{i}: a_{i} \in F, \lambda_{i} \in I\right\}
$$

In this case, we write $K=F[I]$. We define the degree of $\left.K\right|_{p} F$, notation $[K: F]$, by the following $[K: F]:=\infty$ or $[K: F]:=\max \left\{n:\right.$ the set $\left\{1, \lambda, \lambda^{2}, \ldots, \lambda^{n}\right\}$ is linearly independent for all $\left.\lambda \in K\right\}$. rem1
Remark 3.6.6. There are these immediate consequences of the above definitions:
$a-$ If $I \subseteq K$ is linearly independent and $J \subseteq I$ then $J$ is also linearly independent.
$b$ - An element $\alpha \in K$ is $F$-algebraic if and only if $\left\{\alpha^{k}: k \in \mathbb{N}\right\}$ is $F$-linearly dependent.
$c$ - If $[K: F]<\infty$ then all $\alpha \in K$ is $F$-algebraic.
$d$ - Let $F$ be a superfield and $p \in F[X]$ an irreducible polynomial, say $p(X)=a_{0}+a_{1} X+\ldots+$ $a_{n} X^{n-1}+X^{n}$. Then $\left\{\overline{1}, \bar{X}, \ldots, \bar{X}^{n-1}\right\}$ is a $F$-basis of $F(p)$.

Now, let $\left.K\right|_{p} F$ be a proto extension and $\gamma \in K$ algebraic. Then there exist an irreducible polynomial $f(X)$ such that $0 \in f(\gamma, K)$. Let $\operatorname{Irr}_{F}(\gamma, K)$ be the minimum degree irreducible polynomial $f(X)$ such that $0 \in f(\gamma, K)$. Let $F[\gamma, K] \subseteq K$ be the set

$$
F[\gamma, K]:=\bigcup_{f \in F[X]} e v(f, \gamma, K) \subseteq K
$$

and $I_{\gamma, K} \subseteq F[\gamma, K]$ the set

$$
I_{\gamma, K}:=\bigcup_{f \in\left\langle\operatorname{Irr}_{F}(\gamma, K)\right\rangle} e v(f, \gamma, K) \subseteq K
$$

Note that for all $g \in F[X]$ and all $a_{0}, \ldots, a_{n} \in F$, applying the "Newton's binom formula" we get

$$
e v\left(g,\left(a_{0}+a_{1} \gamma+a_{2} \gamma^{2}+\ldots+a_{n-1} \gamma^{n-1}+a_{n} \gamma^{n}\right), K\right) \subseteq F[\gamma, K]
$$

## Remark 3.6.7.

$i$ - If $K \mid F$ is a field extension then our $F[\gamma, K]$ coincide with the usual simple extension $F(\gamma)$.
ii - If $K \mid F$ is a superfield extension and $\gamma \in K$, then $F[\gamma, K]$ depends on the choice of $K$. For example, consider $H_{3} \mid H_{1}$ and $H_{5} \mid H_{1}$ and the element $2 \in H_{3}$ (and of course, in $H_{5}$ ). Then

$$
\begin{aligned}
& H_{2}\left[2, H_{3}\right]=\bigcup_{f \in H_{2}[X]} e v\left(f, \gamma, H_{3}\right)=H_{3}, \\
& H_{2}\left[2, H_{5}\right]=\bigcup_{f \in H_{2}[X]} e v\left(f, \gamma, H_{5}\right)=H_{5},
\end{aligned}
$$

and then, $H_{2}\left[2, H_{3}\right] \neq H_{2}\left[2, H_{5}\right]$.
iii - For a proto extension $\left.K\right|_{p} F$ the set $F[\gamma, K]$ may not be a superfield! Let $F=H_{2}, K=\mathbb{R}$ and $\gamma=2$. Then

$$
H_{2}[2, \mathbb{R}]=2 \mathbb{Z}
$$

which is not a superfield.
At this point, our goal is to obtain an appropriate notion for simple extensions of superfields. In other words, given a full extension $\left.K\right|_{f} F$ and $\alpha \in K$ algebraic, it is highly desirable to obtain a superfield $F(\alpha)$ that:

1. $F \cup\{\alpha\} \subseteq F(\alpha)$;
2. $F(\alpha)$ is the minimal superfield (with respect to inclusion) satisfying (1);
3. $F(\alpha)$ is "computable" in some way (or saying it in a more realistic manner, we want that $F(\alpha) \cong F(p)$ with $\left.p(X)=\operatorname{Irr}_{F}(\alpha)\right)^{6}$.

For general superfields there are some obstacles to achieve this goal. The very first one is the fact that $R[X]$ is not full in general. However, we have an interesting property valid for all $a, b \in R[X]:$

$$
a(1+X)=a+a X \text { and }(a+b) X=a X+b X
$$

This property is the inspiration for the following definition.
almostfull
Definition 3.6.8. Let $\left.K\right|_{p} F$ be a proto superfield extension and $\gamma \in K$. Suppose that $K$ is $F$ generated by $\left\{1, \gamma^{2}, \ldots, \gamma^{n}\right\}$. We say that $K$ is $F$-almost full relative to $\gamma$ (or just almost full) if for all $a, b, c \in F$, and all $p, q, r \in \mathbb{N}$ distinct

$$
\left(a \gamma^{p}+b \gamma^{q}+c \gamma^{r}\right) \gamma=a \gamma^{p+1}+b \gamma^{q+1}+c \gamma^{r+1} .
$$

Here are some immediate consequences of Definition 3.6.8.
lemfator3
Lemma 3.6.9. Let $\left.K\right|_{f} F$ be a full extension $F$-almost full relative to $\gamma$ and let $A=a_{0}+a_{1} \gamma+$ $a_{2}^{2}+\ldots+a_{n} \gamma^{n}$. Then:

$$
i \text { - For all } b, c \in F,(b+c \gamma) A=b A+c \gamma A \text {. }
$$

[^11]ii - For all $b_{0}, \ldots, b_{m} \in F$,
\[

$$
\begin{gathered}
\left(b_{0}+b_{1} \gamma+\ldots+b_{j} \gamma^{j}+b_{j+1} \gamma^{j+1}+\ldots+b_{m} \gamma^{m}\right) A= \\
\left(b_{0}+b_{1} \gamma+\ldots+b_{j} \gamma^{j}\right) A+\left(b_{j+1} \gamma^{j+1}+\ldots+b_{m} \gamma^{m}\right) A .
\end{gathered}
$$
\]

In particular, if $d \in F, B \subseteq K$ with $B=b_{0}+b_{1} \gamma+b_{2} \gamma^{2}+\ldots+b_{m} \gamma^{m}$ and $r>m$, then

$$
\left(B+d \gamma^{r}\right) A=A B+d \gamma^{r} A
$$

Proof. Similar to Lemma 3.4.6.
Lemma 3.6.10. Let $\left.K\right|_{f} F$ be a full extension $F$-almost full relative to $\gamma$. Then:
$i-K=F[\gamma, K]$;
ii - If $\left.K\right|_{f} F$ and $\left.L\right|_{f} K$ are almost full then $L \mid F$ is almost full;
iii - If $\left.L\right|_{f} F$ is another full extension and $\pi: K \rightarrow L$ is a full surjective morphism, then $\left.L\right|_{f} F$ is $F$-almost full relative to $\pi(\gamma)$;
$i v-$ For all $a_{0}, \ldots, a_{n}, b_{0}, \ldots, b_{n} \in F$,

$$
\begin{aligned}
& \left(a_{0}+a_{1} \gamma+a_{2} \gamma^{2}+\ldots+a_{n-1} \gamma^{n-1}+a_{n} \gamma^{n}\right)\left(b_{0}+b_{1} \gamma+b_{2} \gamma^{2}+\ldots+b_{n-1} \gamma^{n-1}+b_{n} \gamma^{n}\right) \subseteq \\
& a_{0} b_{0}+\left(\sum_{j=0}^{1} a_{j} b_{1-j}\right) \gamma+\ldots+\left(\sum_{j=0}^{2 n-1} a_{j} b_{(2 n-1)-j}\right) \gamma^{2 n-1}+\left(\sum_{j=0}^{2 n} a_{j} b_{1-j}\right) \gamma^{2 n}
\end{aligned}
$$

with the convention $a_{j}=b_{j}=0$ if $j>n$.
Let $\left.K\right|_{f} F$ be a full extension and $\alpha \in K$ algebraic over $F$. Our aim is to provide an almost full algebraic extension $\left.F(\alpha)\right|_{f} F$ containing $F$ and $\alpha$. The key to that is to find a way to describe algebraic elements of $K$. Here we have a first result in this direction.

Theorem 3.6.11 (Almost Full Newton's Binom). Let $K \mid F$ be an almost full superfield extension $F$-generated by $\left\{1, \gamma, \ldots, \gamma^{n}\right\}, \gamma \in K$. Then for all $a, b \in F$,

$$
(a+b \gamma)^{n}=\sum_{j=0}^{n}\binom{n}{j} a^{j}(b \gamma)^{n-j}
$$

Proof. By induction is enough to prove the case $n=2$. We have

$$
\begin{aligned}
(a+b \gamma)^{2} & :=(a+b \gamma)(a+b \gamma) \stackrel{\sqrt[3.6 .9]{=}}{=} a(a+b \gamma)+b \gamma(a+b \gamma)=a^{2}+a b \gamma+b \gamma a+(b \gamma)^{2} \\
& =a^{2}+a b \gamma+a b \gamma+(b \gamma)^{2}=a^{2}+2 a b \gamma+(b \gamma)^{2}:=\sum_{j=0}^{2}\binom{n}{j} a^{j}(b \gamma)^{n-j}
\end{aligned}
$$

In the sequence, we have a key result, which states that our "best candidate for simple extension", $F(p)$, is an full algebraic and almost full extension of $F$.

Theorem 3.6.12. Let $F$ be a superfield and $p \in F[X]$ be an irreducible polynomial. Then $\left.F(p)\right|_{f} F$ is an algebraic extension. Moreover, if $\operatorname{deg} p=n+1$ then $[F(p): F] \leq n+1$.

Proof. Remember that $F(p)$ is generated by $\left\{1, \gamma, \ldots, \gamma^{n}\right\}$ with $\gamma=\bar{X}, n \in \mathbb{N}$. Also, we can consider $n$ as the minimal integer such that there exist $a_{0}, \ldots, a_{n+1}$ with

$$
0 \in a_{0}+a_{1} \gamma+\ldots+a_{n+1} \gamma^{n+1}
$$

Now let $b_{0}+b_{1} \gamma+\ldots+b_{n} \gamma^{n} \in F(p)^{*}$. Since $x \cdot y \neq \emptyset$ for all $x, y \in F(P)$, for all $k=0, \ldots, n$, there exist

$$
d_{k 0}+d_{k 1} \gamma+\ldots+d_{k n} \gamma^{n} \in\left(b_{0}+b_{1} \gamma+\ldots+b_{n} \gamma^{n}\right)^{k}
$$

for suitable $d_{i j} \in F$. Writing this in matrix notation, we have

$$
D\left(\begin{array}{c}
1 \\
\gamma \\
\vdots \\
\gamma^{n}
\end{array}\right) \subseteq\left(\begin{array}{c}
\left(b_{0}+b_{1} \gamma+\ldots+b_{n} \gamma^{n}\right)^{0} \\
\left(b_{0}+b_{1} \gamma+\ldots+b_{n} \gamma^{n}\right)^{1} \\
\vdots \\
\left(b_{0}+b_{1} \gamma+\ldots+b_{n} \gamma^{n}\right)^{n}
\end{array}\right)
$$

with

$$
D=\left(\begin{array}{cccc}
d_{00} & d_{01} & \ldots & d_{1 n} \\
d_{10} & d_{11} & \ldots & d_{1 n} \\
\vdots & \vdots & \ddots & \vdots \\
d_{n 0} & d_{n 1} & \ldots & d_{n n}
\end{array}\right)
$$

The fact that $F(p)$ is almost full enable us to scale the matrix $D$, saying

$$
D_{\text {scaled }}=\left(\begin{array}{cccc}
e_{00} & e_{01} & \ldots & e_{1 n} \\
0 & e_{11} & \ldots & e_{1 n} \\
\vdots & \vdots & \ddots & \vdots \\
0 & 0 & \ldots & e_{n n}
\end{array}\right)
$$

and getting

$$
\left(\begin{array}{cccc}
e_{00} & e_{01} & \ldots & e_{1 n}  \tag{*}\\
0 & e_{11} & \ldots & e_{1 n} \\
\vdots & \vdots & \ddots & \vdots \\
0 & 0 & \ldots & e_{n n}
\end{array}\right)\left(\begin{array}{c}
1 \\
\gamma \\
\vdots \\
\gamma^{n}
\end{array}\right) \in\left(\begin{array}{c}
\sum_{j=0}^{n} g_{0 j}\left(b_{0}+b_{1} \gamma+\ldots+b_{n} \gamma^{n}\right)^{j} \\
\sum_{j=0}^{n} g_{1 j}\left(b_{0}+b_{1} \gamma+\ldots+b_{n} \gamma^{n}\right)^{j} \\
\vdots \\
\sum_{j=0}^{n} g_{n j}\left(b_{0}+b_{1} \gamma+\ldots+b_{n} \gamma^{n}\right)^{j}
\end{array}\right)
$$

for suitable $g_{i j} \in F$.

If $D_{\text {scaled }}$ is not invertible then $0 \in \operatorname{det}\left(D_{\text {scaled }}\right)=e_{11} e_{22} \ldots e_{n n}$ and then, $e_{i i}=0$ for some $i \in\{1, \ldots, n\}$ (see Lemma 3.3.11), which imply (by the very scalation process) that there exist a row $i$ with $L_{i}$ being a linear combination of the others. Suppose without loss of generality that

$$
L_{n+1} \cap\left[\left(\sum_{j=1}^{r_{1}} \lambda_{j 1}\right) L_{1}+\ldots+\left(\sum_{j=1}^{r_{n}} \lambda_{j n}\right) L_{n}\right] \neq \emptyset .
$$

This means
$0 \in z_{0}+z_{1}\left(b_{0}+b_{1} \gamma+\ldots+b_{n} \gamma^{n}\right)^{1}+\ldots+z_{n}\left(b_{0}+b_{1} \gamma+\ldots+b_{n} \gamma^{n}\right)^{n-1}-\left(b_{0}+b_{1} \gamma+\ldots+b_{n} \gamma^{n}\right)^{n+1}$,
for suitable $z_{0}, \ldots, z_{n} \in F$, and then, for $f(X)=z_{0}+z_{1} X+\ldots+z_{n} X^{n}-X^{n+1}$, we have

$$
0 \in f\left(b_{0}+b_{1} \gamma+\ldots+b_{n} \gamma^{n}\right)
$$

which means $b_{0}+b_{1} \gamma+\ldots+b_{n} \gamma^{n}$ is algebraic. If $D_{\text {scaled }}$ is invertible, since $\left.F(p)\right|_{f} F$ is almost full we get

$$
\left(\begin{array}{c}
1 \\
\gamma \\
\vdots \\
\gamma^{n}
\end{array}\right) \in D_{\text {scaled }}^{-1}\left[D_{\text {scaled }}\left(\begin{array}{c}
1 \\
\gamma \\
\vdots \\
\gamma^{n}
\end{array}\right)\right]
$$

After multiplying the equation $(*)$ by $D_{\text {scaled }}^{-1}$ we arrive at a system

$$
\left(\begin{array}{c}
1 \\
\gamma \\
\vdots \\
\gamma^{n}
\end{array}\right) \in D_{\text {scaled }}^{-1}\left[D_{\text {scaled }}\left(\begin{array}{c}
1 \\
\gamma \\
\vdots \\
\gamma^{n}
\end{array}\right)\right] \subseteq D_{\text {scaled }}^{-1}\left(\begin{array}{c}
\sum_{j=0}^{n} g_{0 j}\left(b_{0}+b_{1} \gamma+\ldots+b_{n} \gamma^{n}\right)^{j} \\
\sum_{j=0}^{n} g_{1 j}\left(b_{0}+b_{1} \gamma+\ldots+b_{n} \gamma^{n}\right)^{j} \\
\vdots \\
\sum_{j=0}^{n} g_{n j}\left(b_{0}+b_{1} \gamma+\ldots+b_{n} \gamma^{n}\right)^{j}
\end{array}\right)
$$

then our situation is

$$
\left(\begin{array}{c}
1  \tag{**}\\
\gamma \\
\vdots \\
\gamma^{n}
\end{array}\right) \in D_{\text {Scaled }}^{-1}\left(\begin{array}{c}
\sum_{j=0}^{n} g_{0 j}\left(b_{0}+b_{1} \gamma+\ldots+b_{n} \gamma^{n}\right)^{j} \\
\sum_{j=0}^{n} g_{1 j}\left(b_{0}+b_{1} \gamma+\ldots+b_{n} \gamma^{n}\right)^{j} \\
\vdots \\
\sum_{j=0}^{n} g_{n j}\left(b_{0}+b_{1} \gamma+\ldots+b_{n} \gamma^{n}\right)^{j}
\end{array}\right)
$$

Let $D_{\text {Scaled }}^{-1}=\left(h_{i j}\right)$. From $(* *)$, after calculating the matrix product we get (remember the almost fullness)

$$
\left\{\begin{array}{l}
\gamma^{0} \in \sum_{j=0}^{n} g_{0 j} h_{0 j}\left(b_{0}+b_{1} \gamma+\ldots+b_{n} \gamma^{n}\right)^{j} \\
\gamma_{1} \in \sum_{j=0}^{n} g_{1 j} h_{1 j}\left(b_{0}+b_{1} \gamma+\ldots+b_{n} \gamma^{n}\right)^{j} \\
\vdots \\
\gamma^{n} \in \sum_{j=0}^{n} g_{n j} h_{n j}\left(b_{0}+b_{1} \gamma+\ldots+b_{n} \gamma^{n}\right)^{j}
\end{array}\right.
$$

which imply

$$
\left\{\begin{array}{l}
a_{0} \gamma^{0} \in \sum_{j=0}^{n} a_{0}\left(g_{0 j} h_{0 j}\right)\left(b_{0}+b_{1} \gamma+\ldots+b_{n} \gamma^{n}\right)^{j} \\
a_{1} \gamma_{1} \in \sum_{j=0}^{n} a_{1}\left(g_{1 j} h_{1 j}\right)\left(b_{0}+b_{1} \gamma+\ldots+b_{n} \gamma^{n}\right)^{j} \\
\vdots \\
a_{n} \gamma^{n} \in \sum_{j=0}^{n} a_{n}\left(g_{n j} h_{n j}\right)\left(b_{0}+b_{1} \gamma+\ldots+b_{n} \gamma^{n}\right)^{j}
\end{array}\right.
$$

Then

$$
a_{1} a_{n} \gamma^{n+1} \subseteq\left(\sum_{j=0}^{n} a_{1}\left(g_{1 j} h_{1 j}\right)\left(b_{0}+b_{1} \gamma+\ldots+b_{n} \gamma^{n}\right)^{j}\right)\left(\sum_{j=0}^{n} a_{n}\left(g_{n j} h_{n j}\right)\left(b_{0}+b_{1} \gamma+\ldots+b_{n} \gamma^{n}\right)^{j}\right)
$$

and

$$
\begin{aligned}
& 0 \in a_{0}+a_{1} \gamma+\ldots+a_{n+1} \gamma^{n+1} \subseteq \\
& \sum_{p=0}^{n}\left[\sum_{j=0}^{n} a_{p}\left(g_{p j} h_{p j}\right)\left(b_{0}+b_{1} \gamma+\ldots+b_{n} \gamma^{n}\right)^{j}\right]+ \\
& \left(\sum_{j=0}^{n} a_{1}\left(g_{1 j} h_{1 j}\right)\left(b_{0}+b_{1} \gamma+\ldots+b_{n} \gamma^{n}\right)^{j}\right)\left(\sum_{j=0}^{n} a_{n}\left(g_{n j} h_{n j}\right)\left(b_{0}+b_{1} \gamma+\ldots+b_{n} \gamma^{n}\right)^{j}\right) .
\end{aligned}
$$

Now, thinking with polynomials, we have

$$
\begin{aligned}
A(X) & :=\sum_{p=0}^{n}\left[\sum_{j=0}^{n} a_{p}\left(g_{p j} h_{p j}\right) X^{j}\right]+\left(\sum_{j=0}^{n} a_{1}\left(g_{1 j} h_{1 j}\right) X^{j}\right)\left(\sum_{j=0}^{n} a_{n}\left(g_{n j} h_{n j}\right) X^{j}\right)= \\
& {\left[\sum_{p=0}^{n} \sum_{j=0}^{n} a_{p}\left(g_{p j} h_{p j}\right)\right] X^{j}+\left(\sum_{j=0}^{n} a_{1}\left(g_{1 j} h_{1 j}\right) X^{j}\right)\left(\sum_{j=0}^{n} a_{n}\left(g_{n j} h_{n j}\right) X^{j}\right)=} \\
& =P(X)+S(X) T(X),
\end{aligned}
$$

with

$$
\begin{aligned}
& P(X)=\left[\sum_{p=0}^{n} \sum_{j=0}^{n} a_{p}\left(g_{p j} h_{p j}\right)\right] X^{j} \\
& S(X)=\sum_{j=0}^{n} a_{1}\left(g_{1 j} h_{1 j}\right) X^{j} \\
& T(X)=\sum_{j=0}^{n} a_{n}\left(g_{n j} h_{n j}\right) X^{j}
\end{aligned}
$$

Then

$$
0 \in e v\left(A(X), b_{0}+b_{1} \gamma+\ldots+b_{n} \gamma^{n}\right) ;
$$

which means that there exists at least a polynomial $f(X) \in A(X)=P(X)+S(X) T(X)$ with

$$
0 \in f\left(b_{0}+b_{1} \gamma+\ldots+b_{n} \gamma^{n}\right)
$$

Then $b_{0}+b_{1} \gamma+\ldots+b_{n} \gamma^{n}$ is algebraic. Of course, this also imply that $[F(p): F] \leq n+1$.
Keeping on hands the Theorem 3.6.12, we work in order to legitimate $F(p)$ as the simple extension of $F$ by $\alpha$. But before we do that, lets make some considerations about general almost full extensions.

The proof of Theorem 3.6 .12 strongly rely in the fact that $a_{0}+a_{1} \omega+\ldots+a_{n-1} \omega^{n-1}$ is unitary. It is a special property of $F(p)$, and is not necessarily valid for a general almost full full extension.

For an almost full extension $\left.K\right|_{f} F$ denote

$$
\operatorname{Alg}(K, F)=\{\alpha \in K: \alpha \text { is algebraic over } F\} .
$$

We do not know if $\operatorname{Alg}(K, F)$ is a superfield in general. The difficult here is that despite the fact
that Theorem 3.6 .12 is still available, we cannot use it to conclude that all elements in $\alpha \beta$ and $\alpha+\beta$ are algebraic if $\alpha$ and $\beta$ are algebraic.

It is time to define a notion of simple extension.
Definition 3.6.13 (Simple Extension). Let $\left.K\right|_{f} F$ be a full extension and $\alpha \in K$ algebraic. We define the simple extension $F(\alpha, K)$ by

$$
F(\alpha, K):=\bigcap\left\{L:\left.L\right|_{f} F \text { is full and } F[\alpha] \subseteq L\right\}
$$

Note that we have a full extension $\left.F(\alpha, K)\right|_{f} F$. If $\lambda_{1}, \ldots, \lambda_{n} \in K$ are algebraic, we define

$$
F\left(\lambda_{1}, \ldots, \lambda_{n}, K\right):=F\left(\lambda_{1}, \ldots, \lambda_{n-1}, K\right)\left(\lambda_{n}, K\right) .
$$

By Theorem 3.6.3 we can simply write $F(\alpha)$ to indicate $F(\alpha, K)$

## Theorem 3.6.14.

$i$ - Let $\left.K\right|_{f} F$ be a full extension with $\alpha \in K$ algebraic. Let $p(X)=\operatorname{Irr}_{F}(\alpha, K)$. Then $F(\alpha) \cong$ $F(p)$.
ii - Let $\left.K\right|_{f} F$ be a full extension and $\alpha, \beta \in K$ algebraic such that $\left.F(\alpha)(\beta)\right|_{f} F(\alpha)$ and $\left.F(\beta)(\alpha)\right|_{f} F(\beta)$ are almost full extensions relative to $\alpha$ and $\beta$ respectively. Then

$$
F(\alpha)(\beta) \cong F(\beta)(\alpha)
$$

iii - Let $\left.K\right|_{f} F$ be a full extension. For all $\alpha_{1}, \ldots, \alpha_{n} \in K$ and all $\sigma \in S_{n}$ we have

$$
F\left(\alpha_{1}, \ldots, \alpha_{n}\right) \cong F\left(\alpha_{\sigma(1)}, \ldots, \alpha_{\sigma(n)}\right)
$$

Proof.
i - We have that $\left.F(p)\right|_{f} F$ is a full extension containing $F[\alpha, K]$ (see Theorem 3.5.8), so $F(\alpha) \subseteq$ $F(p)$. Moreover, $F(p)$ is generated by $\left\{1, \alpha, \ldots, \alpha^{n-1}\right\}$, where $n=\operatorname{deg}(p)$. Then $F[\alpha]=$ $F\left[\left\{1, \alpha, \ldots, \alpha^{n-1}\right\}\right]$ already is a superfield and

$$
F(p) \cong F\left[\left\{1, \alpha, \ldots, \alpha^{n-1}\right\}\right]=F(\alpha) .
$$

ii - By construction, $e v(p, \alpha, F(\alpha)[X]) \subseteq F(\beta)(\alpha)$ for all $p \in F(\alpha)[X]$. Then $F(\alpha)(\beta) \subseteq F(\beta)(\alpha)$. Reverting the argument we conclude $F(\beta)(\alpha) \subseteq F(\alpha)(\beta)$.
iii - Just use previous item and induction.

Corollary 3.6.15. Let $\left.K\right|_{f} F$ be a full extension with $\alpha \in K$ algebraic and $\operatorname{deg}\left(\operatorname{Irr}_{F}(\alpha)\right)=n$. Then

$$
F(\alpha) \cong\left\{a_{0}+a_{1} \alpha+\ldots+a_{n} \alpha^{n}: a_{0}, \ldots, a_{n} \in F\right\},
$$

with operations in the set on the right inherited from $F[X]$.
Of course, deal with $F(p)$ is much easier to deal with the general expression

$$
\left.\bigcap_{\{L}:\left.L\right|_{f} F \text { is full and } F[\alpha] \subseteq L\right\}
$$

in the sense of make calculations. But the task of determining $F(p)$ "by hand" was already difficult in the field case. In the superfield case this difficult is accentuate, even for low degree polynomials.

Example 3.6.16 (Quadratic Extensions of $H_{3}$ ). Of course, the only irreducible polynomial of degree 2 over $H_{2}$ is $f(X)=X^{2}+2$. We want to describe some possibilities for $H_{3}(\sqrt{2}, K)$ (even in the case of non full extensions).

We first use Theorem 3.5.8. Let $\operatorname{Irr}_{H_{3}}(\sqrt{2})=p(X)=X^{2}+2$ and consider $K=H_{3}(p)$. Lets look closely at the operations on $K$. Denote an element in $K$ by $[f] \in K, f \in H_{3}[X]$. We have

$$
K=\{[0],[1],[2],[X],[2 X],[1+X],[2+X],[1+2 X],[2+2 X]\} .
$$

By definition, for $[f],[g] \in K$ we have

$$
[f]+[g]:=\{[h]: h \in f+g\} \text { and }[f] \cdot[g]:=\{[h]: h \in f g\} .
$$

With these rules is easy to show that $K \mid H_{3}$ is an algebraic full extension (for example, $[1+X]$ is a root of $\left.f(X)=X^{2}+1\right)$. In fact, $K=H_{3}(\sqrt{2})$. Moreover $K$ is not a hyperfield because

$$
([1+X])([1+X])=\dot{K} .
$$

Now let $L=H_{3} \times_{h} H_{5}$. Note that $|L|=(3-1)(5-1)+1=2 \cdot 4+1=9$. Moreover, we have a morphism $i: H_{3} \hookrightarrow H_{5}$ given by the rule $i(x)=\left(1, x^{2}\right)$. Denoting $\omega=(1,2)$, we have

$$
\omega^{2}=(1,2)^{2}=(1,2) \cdot(1,2)=\left(1,2^{2}\right)=(1,4)=i(2) .
$$

More explicitly, doing the following identifications

$$
\begin{aligned}
(1,1) \mapsto 1, & (2,1) \mapsto a, \\
(1,2) \mapsto \omega, & (2,2) \mapsto b, \\
(1,3) \mapsto 2 \omega, & (2,3) \mapsto c, \\
(1,4) \mapsto 2, & (2,4) \mapsto d,
\end{aligned}
$$

we have that

$$
L \cong\{0,1,2, \omega, 2 \omega, a, b, c, d\}
$$

with the following table of operations:

| + | $\omega$ | $2 \omega$ | $a$ | $b$ | $c$ | $d$ |
| :---: | :---: | :---: | :---: | :---: | :---: | :---: |
| 1 | $\{1, \omega, a, b\}$ | $\{1,2 \omega, a, c\}$ | $K \backslash\{0\}$ | $\{1, \omega, a, b\}$ | $\{1,2 \omega, a, c\}$ | $\{1,2, a, d\}$ |
| 2 | $\{2, \omega, b, d\}$ | $\{2,2 \omega, c, d\}$ | $\{1,2, a, d\}$ | $\{2, \omega, b, d\}$ | $\{2,2 \omega, c, d\}$ | $K \backslash\{0\}$ |
| $\omega$ | $K$ | $\{\omega, 2 \omega, b, c\}$ | $\{1, \omega, a, b\}$ | $K \backslash\{0\}$ | $\{\omega, 2 \omega, b, c\}$ | $\{2, \omega, b, d\}$ |
| $2 \omega$ | $\{\omega, 2 \omega, b, c\}$ | $K$ | $\{1,2 \omega, a, c\}$ | $\{\omega, 2 \omega, b, c\}$ | $K \backslash\{0\}$ | $\{2,2 \omega, c, d\}$ |
| $a$ | $\{1, \omega, a, b\}$ | $\{1,2 \omega, a, c\}$ | $K$ | $\{1, \omega, a, b\}$ | $\{1,2 \omega, a, c\}$ | $\{1,2, a, d\}$ |
| $b$ | $K \backslash\{0\}$ | $\{\omega, 2 \omega, b, c\}$ | $\{1, \omega, a, b\}$ | $K$ | $\{\omega, 2 \omega, b, c\}$ | $\{2, \omega, b, d\}$ |
| $c$ | $\{\omega, 2 \omega, b, c\}$ | $K \backslash\{0\}$ | $\{1,2 \omega, a, c\}$ | $\{\omega, 2 \omega, b, c\}$ | $K$ | $\{2,2 \omega, c, d\}$ |
| $d$ | $\{2, \omega, b, d\}$ | $\{2,2 \omega, c, d\}$ | $\{1,2, a, d\}$ | $\{2, \omega, b, d\}$ | $\{2,2 \omega, c, d\}$ | $K$ |


| $\cdot$ | 2 | $\omega$ | $2 \omega$ | $a$ | $b$ | $c$ | $d$ |
| :---: | :---: | :---: | :---: | :---: | :---: | :---: | :---: |
| 2 | 1 | $2 \omega$ | $\omega$ | $d$ | $c$ | $b$ | $a$ |
| $\omega$ | $2 \omega$ | 2 | 1 | $b$ | $d$ | $a$ | $c$ |
| $2 \omega$ | $\omega$ | 1 | 2 | $c$ | $a$ | $d$ | $b$ |
| $a$ | $d$ | $b$ | $c$ | 1 | $\omega$ | $2 \omega$ | 2 |
| $b$ | $c$ | $d$ | $a$ | $\omega$ | 2 | 1 | $2 \omega$ |
| $c$ | $b$ | $a$ | $d$ | $2 \omega$ | 1 | 2 | $\omega$ |
| $d$ | $a$ | $c$ | $b$ | 2 | $2 \omega$ | $\omega$ | 1 |

and of course, $1+1=2+2=L, 0+x=\{x\}, 1 \cdot x=x$ and $0 \cdot x=0$ for all $x \in L$. With these calculations we immediately have that $L$ is an algebraic extension of $H_{3}$.

Now Let $q$ be an odd prime integer greater than 3. The same calculations (with $\omega=(1,2)$ ) proves that $H_{3} \times_{h} H_{q}$ is another algebraic extension of $H_{3}$. Of course, we clearly have $H_{3} \times{ }_{h} H_{5} \nsubseteq H_{3} \times_{h} H_{q}$ for $g \geq 7$. And since all these $H_{3} \times{ }_{h} H_{q}$ are hyperfields and $K$ is a superfield that is not a hyperfield we have $K \nexists H_{3} \times{ }_{h} H_{q}$ for all prime $q \geq 5$. Conclusion: we have infinite non isomorphic algebraic (and non full) hyperfield extensions of $H_{3}$.

### 3.7 Algebraic Closure

As expected, there are some generalizations to the classic notion of algebraic closure for fields.
Definition 3.7.1 (Algebraic Closures). Let $F$ and $K$ be superfields.
$i$ - We say that $K$ is a proto algebraic closure of $F$ if $K$ is algebraically closed and $\left.K\right|_{p} F$ is algebraic.
ii - We say that $K$ is an algebraic closure of $F$ if $K$ is algebraically closed and $K \mid F$ is algebraic.
iii - We say that $K$ is a full algebraic closure of $F$ if $K$ is algebraically closed and $\left.K\right|_{f} F$ is algebraic.

Of course, all these notions coincide if we choose a field $F$.
lemuinq1
Lemma 3.7.2. Let $F$ be a superfield and $\left.K\right|_{f} F$ be an algebraic extension. If $K$ is a full algebraic closure of $F$ then $\left.K\right|_{f} F$ is a maximal full algebraic extension.

Proof. If $\left.K\right|_{f} F$ is not maximal, there is a nontrivial full algebraic extension $\left.L\right|_{f} K$. In particular, there is a nontrivial simple extension $\left.K(\alpha)\right|_{f} K$, then $K$ is not an algebraic closure.

Here we achieve the main result of this present paper.

Theorem 3.7.3 (Existence of the full Algebraic Closure). Let $F$ be a superfield. Then exists a full superfield extension $\left.K\right|_{f} F$ such that $K$ is algebraically closed (and then, a full algebraic closure of $F)$. Moreover, we can choose $K$ in order that $\left.K\right|_{f} F$ is algebraic.

Proof. Let $F$ be a superfield. Consider the following set

$$
A:=\left\{\omega_{i}^{f}: f \in F[X], \operatorname{deg}(f) \geq 1, i=1, \ldots, \operatorname{deg}(f)\right\}
$$

In other words, for each $f$ of degree greater or equal to 1 , we are choosing elements $\omega_{1}^{f}, \ldots, \omega_{\operatorname{deg}(f)}^{f}$ to represent "some possible roots for $f$ ". For each $a \in F, a$ is the root of $f_{a}(X)=X-a$, and hence there is an element $\omega_{1}^{f_{a}} \in A$. Let

$$
\Omega=\left(\mathcal{P}(A) \backslash \bigcup_{a \in F}\left\{\omega_{1}^{f_{a}}\right\}\right) \cup F .
$$

Then $F \subseteq \Omega$. Now, consider all the possible superfields that can be defined on elements of $\Omega$. Denote the set of all such superfields by $\mathcal{E}$. Since $\mathcal{E} \subseteq \Omega$, it is in fact a set, and since $F \in \mathcal{E}$, it is a non-empty set.

Let $\left.E\right|_{f} F$ be an almost full algebraic extension of $F$-generated by $\left\{1, \gamma, \ldots, \gamma^{n}\right\}$ where $\gamma \in E \backslash F$ is a root of $f$ in $F[X]$. In other words, we have $E=F(\gamma)$. Let $\omega \in \Omega \backslash F$. We can "make the variable change" $\gamma \mapsto \omega$ and choose distinct elements for all elements in $F(\gamma)$ in order to get a field $F(\omega) \cong F(\gamma)$, such that $F \subseteq F(\omega) \subseteq \Omega$.

Then, for all almost full algebraic extension $E_{j} \subseteq \Omega$ obtained by the above process, we can take the set

$$
S=\left\{E_{j}: j \in J\right\} .
$$

We have $F \in S$ and $S$ is partially ordered by inclusion.
Let $T=\left\{E_{k j}: k \in K\right\}$ be a chain in $S$ and

$$
W=\bigcup_{k \in K} E_{k j} .
$$

Since $W$ is an algebraic extension of $F$, we get $W \in S$. By Zorn's Lemma, there exist some maximal element $\bar{F} \in S$. We prove that $\bar{F}$ is an algebraic closure of $F$.

In fact, suppose that exists $f(X) \in F[X]$ such that $f$ has no roots in $\bar{F}[X]$. Then, take $\omega \in \Omega$ such that $\omega \notin \bar{F}$ and $\omega$ is a root of $f(X)$. Consider the field $\bar{F}(\omega)$ as we did above. Then $\bar{F}(\omega)$ is an algebraic extension with $\bar{F} \subsetneq \bar{F}(\omega)$, contradicting the maximality of $\bar{F}$, which complete the proof.

We are surprisingly able to prove the uniqueness of full algebraic closures.
uniq
Theorem 3.7.4 (Uniqueness of the full Algebraic Closure). Let $F$ be a superfield. Let $K_{1}, K_{2}$ be two full algebraic closures of $F$. Then $K_{1} \cong K_{2}$.

To prove Theorem 3.7.4 we need two Lemmas. Let $\left.L\right|_{f} F$ be a full superfield extension and $N$ be another superfield. An $F$-embedding is a full embedding $\iota: L \rightarrow N$ such that $\iota(a)=a, a \in F$. uniq1
Lemma 3.7.5. Let $\left.L\right|_{f} F$ be an algebraic full extension and $\left.N\right|_{f} L$ another algebraic full extension, and $\bar{F}$ some full algebraic closure of $F$. There is a $F$-embedding $i: L \rightarrow \bar{F}$ and once $i$ is picked there exists a $F$-embedding $N \rightarrow \bar{F}$ extending $i$.
Proof. Since a full embedding $i: L \rightarrow \bar{F}$ realizes the full algebraically closed $\bar{F}$ as an algebraic extension of $L$ (and hence as a full algebraic closure of $L$ ), by renaming the base superfield as $L$ it suffices to just prove the first part: any strong algebraic extension admits a full embedding into a specified full algebraic closure.

Let $\Sigma$ to be the set of pairs $(K, i)$ such that $\left.K\right|_{f} F,\left.L\right|_{f} K$ and the inclusion map $i: K \rightarrow \bar{F}$ is a $F$-embedding. Of course, $(F, i) \in \Sigma$, and using the partial order defined by

$$
\left(K_{1}, i_{1}\right) \leq\left(K_{2}, i_{2}\right) \text { iff }\left.K_{2}\right|_{f} K_{1},\left.L\right|_{f} K_{2} \text { and }\left.i_{2}\right|_{k_{1}}=i_{1},
$$

we obtain that every chain has an upper bound (the superfield obtained by directed union). Then we are under the hypothesis of Zorn's Lemma and there exists a maximal element $(N, i) \in \Sigma$.

We just have to show $N=L$. Pick $\alpha \in L$, so $\alpha$ is algebraic over $N$ (as it is algebraic over $F)$. We have $\left.N(\alpha)\right|_{f} N$ and $\left.\bar{F}\right|_{f} N(\alpha)$. In other words, the inclusion map $i: N(\alpha) \rightarrow \bar{F}$ is a full $N$-embedding. By maximality of $N$ we get $N(\alpha)=N$ for all $\alpha \in L$, which imply $N=L$.
uniq2
Lemma 3.7.6. Let $F$ be a superfield and $\bar{F}$ be some full algebraic closure of $F$. If $\phi: \bar{F} \rightarrow \bar{F}$ is a $F$-embedding then $\phi$ is an isomorphism.

Proof. We only need to show that $\phi$ is surjective. Let $\gamma \in \bar{F}$. Then there exist $p(X) \in F[X]$, saying $p(X)=X^{n}+a_{n-1} X^{n-1}+\ldots+a_{1} X+a_{0}$ with $0 \in p(\gamma)$. Since $\phi$ is a $F$-embedding, we have $p^{\phi}(X):=X^{n}+\phi\left(a_{n-1}\right) X^{n-1}+\ldots+\phi\left(a_{1}\right) X+\phi\left(a_{0}\right)=X^{n}+a_{n-1} X^{n-1}+\ldots+a_{1} X+a_{0}=p(X)$.

Then $\phi(\gamma)$ is a root of $p(X)$ because

$$
\begin{aligned}
& 0 \in a_{n} \gamma^{n}+a_{n-1} \gamma^{n-1}+\ldots+a_{1} \gamma+a_{0} \Rightarrow \phi(0) \in \phi\left(a_{n} \gamma^{n}+a_{n-1} \gamma^{n-1}+\ldots+a_{1} \gamma+a_{0}\right) \Rightarrow \\
& 0 \in a_{n} \phi(\gamma)^{n}+a_{n-1} \phi(\gamma)^{n-1}+\ldots+a_{1} \phi(\gamma)+a_{0} .
\end{aligned}
$$

Since $\phi$ is a full embedding, we have a full embedding $\phi(\bar{F}) \hookrightarrow \bar{F}$. Then $\left.\bar{F}\right|_{f} \phi(\bar{F})$. Since $\bar{F}$ is algebraically closed, every non-constant polynomial $p(X) \in F[X]$ has a root $\gamma \in \bar{F}$, and then, a root $\phi(\gamma) \in \phi(\bar{F})$. If $\phi(\bar{F}) \neq \bar{F}$, we have a contradiction with the maximality of $\phi(\bar{F})$ obtained in Lemma 3.7.2.

Proof of Theorem 3.7.4. By Lemma 3.7.5 applied to $L=K_{1}$ and $\bar{F}=K_{2}$ (a full algebraic closed superfield equipped with a structure of algebraic extension of $F$ ), there exists a $F$-embedding $i_{1}: K_{1} \rightarrow K_{2}$. By the very same argument, there also exists a $F$-embedding $i_{2}: K_{2} \rightarrow K_{1}$. Moreover, $i_{1} \circ i_{2}: K_{1} \rightarrow K_{1}$ and $i_{2} \circ i_{1}: K_{2} \rightarrow K_{2}$ are $F$-embeddings. By Lemma 3.7.6, both $i_{1} \circ i_{2}$ and $i_{2} \circ i_{1}$ are isomorphisms, implying that $i_{1}$ and $i_{2}$ are also isomorphisms.

Example 3.7.7. Lets look at $H_{3}$ again. Consider $L_{1}=H_{3} \times_{h} H_{5}$ and $L_{2}=H_{3} \times_{h} H_{7}$. We do not know precisely the relations between the full algebraic closures $\overline{H_{3}}, \overline{L_{1}}$ and $\overline{L_{2}}$.

Of course, since $L_{1} \mid H_{3}$ and $L_{2} \mid H_{3}$ are algebraic extensions of $H_{3}$, we have that $\overline{L_{1}}$ and $\overline{L_{2}}$ are algebraic closures of $\overline{H_{3}}$. Since $L_{2}$ is an algebraic extension of $L_{1}$, we know that $\overline{L_{2}}$ is an algebraic closure of $L_{1}$. But we do not know if $\overline{H_{3}}, \overline{L_{1}}$ and $\overline{L_{2}}$ are isomorphic (or not).

It is desirable to achieve explicit calculations of $\bar{F}$ for some cases: $F$ (reduced) special hyperfields/groups, in particular $F=\{-1,0,1\}$ and $F$ the special hyperfield/group associated to a Boolean algebra, etc.

### 3.8 Vector Spaces

Since we already have available matrices and polynomials for superrings, a natural extension for the theory is a sort of "vector space" and some linear algebra methods. We start this program here, proceeding in a very similar fashion of Hofmann's and Kunze's Linear Algebra Book ([39).

Throughout this Section, all superfields will be considered associative.
Definition 3.8.1. A (multi) vector space over a superfield $F$ is a tuple ( $V,+, \cdot, 0$ ) such that $(V,+, 0)$ is an abelian multigroup and $: F \times V \rightarrow \mathcal{P}^{*}(V)$ is a function (which image denoted by $\cdot(\lambda, v):=\lambda v)$ satisfying the following properties for all $\lambda, \mu \in F$ and all $v, w \in V$ :

MV0-1v $=\{v\}$ and $0 \cdot v=\{0\}$;
MV1- $\lambda(\mu v)=(\lambda \mu) v$.
Here we adopt the following convention: if $A \subseteq F$ and $v \in V$, we set

$$
A v:=\bigcup\{\lambda v: \lambda \in A\} .
$$

MV2 - $\lambda(v+w) \subseteq \lambda v+\lambda w$;
MV3- $(\lambda+\mu) v \subseteq \lambda v+\mu v$.
The vector space $(V,+, 0)$ is full if the equality holds in MV2 and MV3.
We proceed similarly to the practice used with polynomials and matrices: we omit the word "multi" and just say "vector spaces" over superfields.

Of course, we stick to vectors spaces here but it is available the Definition for "modules", just replacing superfields in Definition 3.8.1 for superring.

Here are some natural examples of vector spaces.

## extvec

Proposition 3.8.2. Let $K \mid F$ be a superfield extension with $K$ and $F$ associative. Then $K$ is a $F$-vector space, which is full iff the extension is full.

Proof. Here $\cdot: F \rightarrow K \rightarrow \mathcal{P}^{*}(K)$ is just the restriction of multiplication to $F$ on the first coordinate. M0 is immediate and M1-M3 are consequences of the axioms of superrings. It is immediate that $K$ is a full vector space iff $\left.K\right|_{f} F$.

Theorem 3.8.3. Let $F^{n}$ be the usual $n$-folded cartesian product $F \times \ldots \times F$. We already know that $F^{n}$ with the induced sum is a multigroup. Now, for $\lambda \in F$ and $v=\left(x_{1}, \ldots, x_{n}\right) \in F^{n}$ define

$$
\lambda v:=\left(\lambda x_{1}, \ldots, \lambda x_{n}\right):=\bigcup\left\{\left(a_{1}, a_{2}, \ldots, a_{n}\right): a_{j} \in \lambda x_{j}, j \geq 1\right\}
$$

Then $\left(F^{n},+, \cdot, 0\right)$ is a vector space. Moreover $F^{n}$ is full iff $F$ is full.
Proof. We already have that $F^{n}$ is commutative a superring. By the very Definition of scalar product we get $1 v=v$. Now let $v, w \in F^{n}, v=\left(x_{1}, \ldots, x_{n}\right), w=\left(y_{1}, \ldots, y_{n}\right)$ and $\lambda, \mu \in F$. We have

$$
\begin{aligned}
(\lambda+\mu) v & :=\left((\lambda+\mu) x_{1}, \ldots,(\lambda+\mu) x_{n}\right) \\
& \subseteq\left(\lambda x_{1}+\mu x_{1}, \ldots, \lambda x_{n}+\mu x_{n}\right) \\
& =\left(\lambda x_{1}, \ldots, \lambda x_{n}\right)+\left(\mu x_{1}, \ldots, \mu x_{n}\right)=\lambda v+\mu w .
\end{aligned}
$$

Similarly we conclude that $(\lambda+\mu) v \subseteq \lambda v+\mu v$.
Then $F^{n}$ is a vector space which is full if $F$ is full.
Now suppose $F^{n}$ full. Then for $\alpha, \lambda, \mu \in F$ and $v=(\alpha, 0, \ldots, 0)$ we have

$$
(\lambda \alpha+\mu \alpha, 0, \ldots, 0)=\lambda v+\lambda v=(\lambda+\mu) v=((\lambda+\mu) \alpha, 0, \ldots, 0)
$$

which means $(\lambda+\mu) \alpha=\lambda \alpha+\mu \alpha$. Similarly we conclude that $\alpha(\lambda+\mu)=\alpha \lambda+\alpha \mu$.
Theorem 3.8.4. Let $F$ be a superfield and $m, n \geq 1$. Then $M_{m \times n}(F)$ is a vector space which is full iff $F$ is full.

Proof. This is consequence of Lemma 3.2.7, identifying $F$ with $M_{1 \times 1}(F)$.

Theorem 3.8.5. Let $F$ be a superfield and $n \geq 1$. Then $F\left[X_{1}, \ldots, X_{n}\right]$ is a vector space which is full iff $F$ is full.

Proof. The argument here is similar to the one in Theorem 3.8.3.
Definition 3.8.6 (Subspace). Let $V$ be a $F$-vector space and $W \subseteq V$. We say that $W$ is a subspace if $0 \in W$ and for all $w_{1}, w_{2} \in W$ and all $\lambda \in F$ we have $w_{1}+w_{2} \subseteq W$ and $\lambda w_{1} \subseteq W$.

Theorem 3.8.7. Let $F$ be a full superfield and consider a system $A x=0, A \in M_{n \times m}(F)$. Then

$$
\operatorname{Sol}[A x=0]:=\left\{v \in M_{n \times 1}(F): 0 \in A v\right\}
$$

is a subspace of $M_{n \times 1}(F)$.
Proof. This is another consequence of Lemma 3.2.7. We need $F$ full in order to conclude that if $0 \in A v$ and $0 \in A w$ then $0 \in A(v+w)=A v+A w$.

Definition 3.8.8 (Spanned Subspace). Let $V$ be a $F$-vector space and $A \subseteq V$. The subspace generated by $A$ is defined by

$$
\langle A\rangle:=\bigcap\{W \subseteq V: W \text { is a subspace and } A \subseteq W\}
$$

Definition 3.8.9 (Linear Combination). Let $V$ be a $F$-vector space, $A \subseteq V$ and $w \in V$. We say that $w$ is a linear combination of elements in $A$ if there exist $\left\{v_{1}, \ldots, v_{n}\right\} \subseteq A$ with

$$
w \in \sum_{j=1}^{r_{1}} \lambda_{j 1} v_{1}+\ldots+\sum_{j=1}^{r_{n}} \lambda_{j n} v_{n}
$$

for some $\lambda_{i j} \in F$. We denote the set of linear combinations of $V$ by

$$
\mathcal{C} \mathcal{L}(A)=\bigcup\left\{\sum_{j=1}^{r_{1}} \lambda_{j 1} v_{1}+\ldots+\sum_{j=1}^{r_{n}} \lambda_{j n} v_{n}:\left\{v_{1}, \ldots, v_{n}\right\} \subseteq A, \lambda_{i j} \in F, r_{1}, \ldots, r_{n} \in \mathbb{N}\right\}
$$

If $V$ is full, then

$$
\mathcal{C} \mathcal{L}(A):=\bigcup\left\{\lambda_{1} v_{1}+\ldots+\lambda_{n} v_{n}: v_{i} \in A, \lambda_{i} \in F, i=1, \ldots, n, n \geq 1\right\}
$$

gen1
Theorem 3.8.10. Let $V$ be a $F$-vector space and $A \subseteq V$. Then $\langle A\rangle=\mathcal{C} \mathcal{L}(A)$.
Proof. We have that $\mathcal{C} \mathcal{L}(A)$ is a subspace, which provide $\langle A\rangle \subseteq \mathcal{C} \mathcal{L}(A)$. If $W \subseteq V$ and $A \subseteq W$, by the very Definition of subspace (and induction) we have $\mathcal{C} \mathcal{L}(A) \subseteq W$, which provide $\mathcal{C} \mathcal{L}(A) \subseteq$ $\langle A\rangle$.

Lemma 3.8.11. Let $V$ be a $F$-vector space and $A, B \subseteq V$. Then

$$
\begin{aligned}
& i-\langle\langle A\rangle\rangle=\langle A\rangle \\
& \text { ii - if } A \subseteq B \text { then }\langle A\rangle \text { is a subspace of }\langle B\rangle \\
& \text { iii - if } A \subseteq B \text { and for all } v \in B, v \in\langle A\rangle \text { then }\langle A\rangle=\langle B\rangle
\end{aligned}
$$

Definition 3.8.12 (Linear Independence). Let $V$ be a $F$-vector space and $A \subseteq V$. We say that $A$ is $F$-linearly independent if for all distinct $v_{1}, \ldots, v_{n} \in A, n \in \mathbb{N}$, the following hold:

$$
\text { If } 0 \in \sum_{j=1}^{r_{1}}\left(\lambda_{j 1} v_{1}\right)+\ldots+\sum_{j=1}^{r_{n}}\left(\lambda_{j n} v_{n}\right) \text { then } 0 \in \sum_{j=1}^{r_{1}} \lambda_{j i} \text { for all } i=1, \ldots, n \text {. }
$$

and $I$ is $F$-linearly dependent if it is not $F$-linearly independent.
Definition 3.8.13 (Base). Let $V$ be a $F$-vector space and $B \subseteq V$. We say that $B$ is a $F$-basis if $B$ is linear independent and $V=\langle B\rangle$.

Definition 3.8.14. We say that a F-vector space is finitely generated if $V=\langle S\rangle$ for some $S \subseteq V$ finite.
basis1
Theorem 3.8.15. Let $F$ be a hyperfield and $V$ be a finitely generated $F$-vector. If $V$ is full then $V$ has a basis.

Proof. Let $F$ be a hyperfield and $V$ be a finitely generated $F$-vector space with $V=\left\langle v_{1}, \ldots, v_{n}\right\rangle$ $\left(v_{1}, \ldots, v_{n} \in V\right)$. If $\left\{v_{1}, \ldots, v_{n}\right\}$ is LI we are done. If not, after a rearrangement of indexes if necessary, we can suppose without loss of generality that $v_{1} \in \mathcal{C} \mathcal{L}\left(\left\{v_{2}, \ldots, v_{n}\right\}\right)$. Then (using Theorem 3.8.10 and Lemma 3.8.11 we have

$$
V=\mathcal{C} \mathcal{L}\left(\left\{v_{1}, \ldots, v_{n}\right\}\right)=\mathcal{C} \mathcal{L}\left(\left\{v_{2}, \ldots, v_{n}\right\}\right) .
$$

If $\left\{v_{2}, \ldots, v_{n}\right\}$ we are done. If not, suppose without loss of generality that $v_{2} \in \mathcal{C} \mathcal{L}\left(\left\{v_{3}, \ldots, v_{n}\right\}\right)$. Then we have

$$
V=\mathcal{C} \mathcal{L}\left(\left\{v_{1}, \ldots, v_{n}\right\}\right)=\mathcal{C} \mathcal{L}\left(\left\{v_{2}, \ldots, v_{n}\right\}\right)=\mathcal{C} \mathcal{L}\left(\left\{v_{3}, \ldots, v_{n}\right\}\right) .
$$

Repeating this process, after a number finite of steps we arrive at a basis $\left\{v_{k}, v_{k+1}, \ldots, v_{n}\right\}$ of $V$ for some $k$ with $1 \leq k \leq n$.

Unfortunately, we do not know if, for general superfields $F$, all basis in a finitely generated $F$ vector spaces has the same dimension. In order to deal with this question, we propose the following concept.
linearly-closed

Definition 3.8.16. Let $F$ be superfield. We say that $F$ is linearly closed if the system $A x=0$ has at least a non trivial solution weak solution for all $A \in M_{n \times m}(F)$ with $m>n$.

Of course, every field is a linearly closed superfield. As we will see later (Theorem 3.8.23), this is also the case for hyperbolic and double distributive hyperfields. The concept of linearly closeness is useful to get the notion of dimension for a subclass of finitely generated $F$-vector spaces.

Definition 3.8.17. Let $F$ be a linearly closed superfield and $V$ be a full $F$-vector space with $V=\langle A\rangle$. We say that $V$ is rigidly generated by $A$ if for all $w \in V$ there exists $v_{i_{1}}, \ldots, v_{i_{n}} \in A$ and $\lambda_{1}, \ldots, \lambda_{n} \in F$ with

$$
\{w\}=\lambda_{1} v_{i_{1}}+\ldots+\lambda_{n} v_{i_{n}} .
$$

Theorem 3.8.18. Let $F$ be a linearly closed superfield and $V$ be a full and finitely generated $F$ vector space with $V=\left\langle v_{1}, \ldots, v_{n}\right\rangle\left(v_{1}, \ldots, v_{n} \in V\right)$. If $V$ is rigidly generated by $\left\{v_{1}, \ldots, v_{n}\right\}$ then every linear independent subset of $V$ has at most $n$ elements.

Proof. We just need to prove that if $S \subseteq V$ and $|S|>n$ then $S$ is linearly dependent.
Let $S$ be such set with $S=\left\{w_{1}, \ldots, w_{m}\right\}, m>n$. Since $V=\left\langle v_{1}, \ldots, v_{n}\right\rangle$, then there exists scalars $a_{i j} \in F$ with

$$
w_{j}=a_{1 j} v_{1}+\ldots+a_{n j} v_{n}, j=1, \ldots, m
$$

Then for all $\lambda_{1}, \ldots, \lambda_{m} \in F$ we get

$$
\lambda_{1} w_{1}+\ldots+\lambda_{m} w_{m}=\sum_{j=1}^{m} \lambda_{j} w_{j}=\sum_{j=1}^{m} \lambda_{j}\left(\sum_{i=1}^{n} a_{i j} v_{i}\right)=\sum_{j=1}^{m} \sum_{i=1}^{n}\left(\lambda_{j} a_{i j}\right) v_{i}=\sum_{i=1}^{n}\left[\sum_{j=1}^{m} \lambda_{j} a_{i j}\right] v_{i}
$$

Then $0 \in \lambda_{1} w_{1}+\ldots+\lambda_{m} w_{m}$ iff

$$
0 \in \sum_{i=1}^{n}\left[\sum_{j=1}^{m} \lambda_{j} a_{i j}\right] v_{i}
$$

providing

$$
0 \in\left(\begin{array}{cccc}
a_{11} & a_{21} & \ldots & a_{m 1} \\
a_{21} & a_{22} & \ldots & a_{m 2} \\
\vdots & \vdots & \ddots & \vdots \\
a_{1 n} & a_{2 n} & \ldots & a_{m n}
\end{array}\right)\left(\begin{array}{c}
\lambda_{1} \\
\lambda_{2} \\
\vdots \\
\lambda_{m}
\end{array}\right)
$$

Let

$$
A=\left(\begin{array}{cccc}
a_{11} & a_{21} & \ldots & a_{m 1} \\
a_{21} & a_{22} & \ldots & a_{m 2} \\
\vdots & \vdots & \ddots & \vdots \\
a_{1 n} & a_{2 n} & \ldots & a_{m n}
\end{array}\right)
$$

Since $F$ is linearly closed, the system $A x=0$ has a weak solution if $m>n$, and we have that $S$ is linear dependent if $m>n$.

Definition 3.8.19. Let $F$ be a linearly closed superfield and $V$ be a full $F$-vector space. We say that $\left\{v_{1}, \ldots, v_{n}\right\}$ is a rigid basis of $V$ if $\left\{v_{1}, \ldots, v_{n}\right\}$ is LI and $V$ is rigidly generated by $\left\{v_{1}, \ldots, v_{n}\right\}$.

Example 3.8.20. For a hyperfield $F$, the $F$-vector spaces $F^{n}, M_{m \times n}(F)$ and $F\left[X_{1}, \ldots, X_{n}\right]$ are all rigidly generated by the analogous canonical basis. In fact, the respective canonical basis is a rigid basis for those spaces.

Theorem 3.8.21. Let $F$ be a linearly closed superfield and $V$ be a $F$-vector space. If $B_{1}$ and $B_{2}$ are rigid basis of $V$ then $\left|B_{1}\right|=\left|B_{2}\right|$.

Proof. Let $B_{1}=\left\{v_{1}, \ldots, v_{n}\right\}$ and $B_{2}=\left\{w_{1}, \ldots, w_{m}\right\}$. Since $V=\left\langle B_{1}\right\rangle$ and $B_{2}$ is linearly independent, by Theorem 3.8.18 we get $m \leq n$. Since $V=\left\langle B_{2}\right\rangle$ and $B_{1}$ is linearly independent, by Theorem 3.8.18 we get $n \leq m$. Then $m=n$.

Definition 3.8.22. Let $F$ be a linearly closed superfield and $V$ be a $F$-vector space finitely generated with a rigid basis. We define the dimension of $V$ by $\operatorname{dim}(V):=|B|$ where $B \subseteq V$ is any rigid basis of $V$.

Of course, it is not clear whether or not a superfield is linearly closed. In the sequence we provide some surprisingly examples, provenient from the structures which we were working until now: hyperbolic hyperfields, which arise naturally in the context of abstract theories of quadratic forms. In particular, there is available the machinery of K-theory for hyperbolic hyperfields ([18]).

Theorem 3.8.23. Every hyperbolic hyperfield is a linearly closed superfield.
Proof. First, remember that every hyperfield is full. Now, let $A \in M_{n \times m}(F)$ with $m>n$. Second, since $F$ is a hyperfield, considering the vector space $F^{n}$, we have a full $F$-vector space such that for all $\lambda \in F$ and all $v \in F^{n}, \lambda v$ is a singleton set.

Keeping this in mind, we will construct a weak solution $b \in M_{m \times 1}(F)$ dividing the proof in some cases.

Case I- $n=1$ (which imply $m \geq 2$ ). In this case, let $A \in M_{1 \times m}(F)$ with

$$
A=\left(\begin{array}{llll}
a_{1} & a_{2} & \ldots & a_{m}
\end{array}\right)
$$

We need to find $x_{1}, \ldots, x_{m} \in F$ (not all zero) such that

$$
0 \in a_{1} x_{1}+a_{2} x_{2}+\ldots+a_{m} x_{m} .
$$

If $a_{i}=0$ for some $i$, just choose $x_{i}=1$ and $x_{j}=0$ for all $j \neq i$. Otherwise, choose $x_{i}=0$ for all $i \geq 3, x_{2}=1$ and $x_{1}=-a_{1}^{-1} a_{2}$. Then

$$
a_{1} x_{1}+a_{2} x_{2}+a_{3} x_{3}+\ldots+a_{m} x_{m}=a_{1} x_{1}+a_{2} x_{2}=a_{1}\left[-a_{1}^{-1} a_{2}\right]+a_{2}=a_{2}-a_{2} \text { with } 0 \in a_{2}-a_{2} .
$$

We further refer to this tactic to find $x_{1}, x_{2}, \ldots, x_{m}$ as "the Case I method" ${ }^{\text {l }}$. ${ }^{\text {ase-I-method }}$

Case II - $n=2$ (which imply $m \geq 3$ ). In this case, let $A \in M_{2 \times m}(F)$ with

$$
A=\left(\begin{array}{cccc}
a_{1} & a_{2} & \ldots & a_{m} \\
b_{1} & b_{2} & \ldots & b_{m}
\end{array}\right)
$$

We need to find $x_{1}, \ldots, x_{m} \in F$ (not all zero) such that

$$
\begin{align*}
& 0 \in a_{1} x_{1}+a_{2} x_{2}+\ldots+a_{m} x_{m} \\
& 0 \in b_{1} x_{1}+b_{2} x_{2}+\ldots+b_{m} x_{m}
\end{align*}
$$

We have some cases here. If $a_{j}=b_{j}=0$ for some $j$, just choose $x_{j}=1$ and $x_{i}=0$ for all $i \neq j$. Now let $0 \in\left\{a_{1}, \ldots, a_{m}, b_{1}, \ldots, b_{m}\right\}$, saying $b_{1}=0$ and $a_{1} \neq 0$. Then we are reduced to

$$
\begin{aligned}
& 0 \in a_{1} x_{1}+a_{2} x_{2}+\ldots+a_{m} x_{m} \\
& 0 \in b_{2} x_{2}+\ldots+b_{m} x_{m}
\end{aligned}
$$

$$
\text { eq-02-sys-case-IT }(3.4)
$$

Now find some $d_{2}, \ldots, d_{m} \in F$ (not all zero) with $0 \in b_{2} d_{2}+\ldots+b_{m} d_{m}$ (as in the Case I method) and choose $d_{1}$ in order to get $d_{1} \in-a_{1}^{-1}\left[a_{2} d_{2}+\ldots+a_{m} d_{m}\right]$. So $\left(d_{1}, \ldots, d_{m}\right)$ is a non trivial weak solution of the reduced system 3.4. Now, let $0 \notin\left\{a_{1}, \ldots, a_{m}, b_{1}, \ldots, b_{m}\right\}$. If $\left\{\left(a_{1}, \ldots, a_{m}\right),\left(b_{1}, \ldots, b_{m}\right)\right\}$ is LD, saying $\left(b_{1}, \ldots, b_{m}\right)=\lambda\left(a_{1}, \ldots, a_{m}\right)$, we just need to choose $d_{1}, d_{2}, \ldots, d_{m} \in F$ not all zero (as in case I) such that $0 \in a_{1} d_{1}+a_{2} d_{2}+\ldots+a_{m} d_{m}$ in order to get a solution of 3.3 for this case.

Now, suppose $0 \notin\left\{a_{1}, \ldots, a_{m}, b_{1}, \ldots, b_{m}\right\}$ and $\left\{\left(a_{1}, \ldots, a_{m}\right),\left(b_{1}, \ldots, b_{m}\right)\right\}$ LI. Since $F$ is a hyperfield, to find a non trivial weak solution of 3.3 is equivalent to find a non trivial weak solution of

$$
\begin{aligned}
& 0 \in x_{1}+a_{1}^{-1} a_{2} x_{2}+\ldots+a_{1}^{-1} a_{m} x_{m} \\
& 0 \in x_{1}+b_{1}^{-1} b_{2} x_{2}+\ldots+b_{1}^{-1} b_{m} x_{m}
\end{aligned}
$$

Then we can suppose without loss of generality that $a_{1}=b_{1}=1$, and we need to find a weak solution of

$$
\begin{aligned}
& 0 \in x_{1}+a_{2} x_{2}+\ldots+a_{m} x_{m} \\
& 0 \in x_{1}+b_{2} x_{2}+\ldots+b_{m} x_{m}
\end{aligned}
$$

Now, consider the set of systems obtained after the elementary operation $L_{2} \leftarrow L_{2}-L_{1}$ :

$$
\begin{align*}
& 0 \in x_{1}+a_{2} x_{2}+\ldots+a_{m} x_{m} \\
& 0 \in(1-1) x_{1}+\left(b_{2}-a_{2}\right) x_{2}+\ldots+\left(b_{m}-a_{m}\right) x_{m}
\end{align*}
$$

As in Case I, let $d_{2}, \ldots, d_{m} \in F$ (not all zero) with

$$
0 \in\left(b_{2}-a_{2}\right) d_{2}+\ldots+\left(b_{m}-a_{m}\right) d_{m} .
$$

In particular, there exist $z \in F$ with

$$
z \in\left(a_{2} d_{2}+\ldots+a_{m} d_{m}\right) \cap\left(b_{2} d_{2}+\ldots+b_{m} d_{m}\right) .
$$

Let $x_{1}=-z$. Then $\left(-z, d_{2}, d_{3}, \ldots, d_{m}\right)$ is a weak solution of both 3.5 and 3.6. completing the proof for Case II. We further refer to this tactic to find $z$ as "the Case II method"l.

Case III - $n=3$ (which imply $m \geq 4$ ). In this case, let $A \in M_{3 \times m}(F)$ with

$$
A=\left(\begin{array}{cccc}
a_{1} & a_{2} & \ldots & a_{m} \\
b_{1} & b_{2} & \ldots & b_{m} \\
c_{1} & c_{2} & \ldots & c_{m}
\end{array}\right)
$$

We need to find $x_{1}, \ldots, x_{m} \in F$ (not all zero) such that

$$
\begin{aligned}
& 0 \in a_{1} x_{1}+a_{2} x_{2}+\ldots+a_{m} x_{m} \\
& 0 \in b_{1} x_{1}+b_{2} x_{2}+\ldots+b_{m} x_{m} \\
& 0 \in c_{1} x_{1}+c_{2} x_{2}+\ldots+c_{m} x_{m}
\end{aligned}
$$

Choosing $x_{j}=0$ for $j \geq 5$, we are reduced to the system

$$
\begin{aligned}
& 0 \in a_{1} x_{1}+a_{2} x_{2}+a_{3} x_{3}+a_{4} x_{4} \\
& 0 \in b_{1} x_{1}+b_{2} x_{2}+b_{3} x_{3}+b_{4} x_{4} \\
& 0 \in c_{1} x_{1}+c_{2} x_{2}+c_{3} x_{3}+c_{4} x_{4}
\end{aligned}
$$

We have some subcases to deal with:
If $a_{j}=b_{j}=c_{j}=0$ for some $j$, just choose $x_{j}=1$ and $x_{i}=0$ for all $i \neq j$.

The second subcase is the one with two elements in $\left\{a_{j}, b_{j}, c_{j}\right\}$ are equal to zero for some $j$. Say for instance that $j=1$ and $a_{1} \neq 0, b_{1}=c_{1}=0$ (the other cases are analogous). Our situation
here is

$$
\begin{aligned}
& 0 \in a_{1} x_{1}+a_{2} x_{2}+a_{3} x_{3}+a_{4} x_{4} \\
& 0 \in b_{2} x_{2}+b_{3} x_{3}+b_{4} x_{4} \\
& 0 \in c_{2} x_{2}+c_{3} x_{3}+c_{4} x_{4}
\end{aligned}
$$

Then we just apply the Case II method 3.8 to get a solution $\left(d_{2}, \ldots, d_{m}\right)$ of the last two equations and for

$$
d_{1} \in-a_{1}^{-1} a_{2} d_{2}-a_{1}^{-1} a_{3} d_{3}-\ldots-a_{1}^{-1} a_{m} d_{3}
$$

we have that $\left(d_{1}, d_{2}, \ldots, d_{m}\right)$ is a solution of 3.7 .

The third subcase is the one with for all $j$, only one element in $\left\{a_{j}, b_{j}, c_{j}\right\}$ is equal to zero. By the pigeon hole principle, one of the sets $\left\{a_{1}, a_{2}, a_{3}, a_{4}\right\},\left\{b_{1}, b_{2}, b_{3}, b_{4}\right\}$ and $\left\{c_{1}, c_{2}, c_{3}, c_{4}\right\}$ has two elements equal to zero. For instance, say that $a_{1}=a_{4}=0, b_{2}=0$ and $c_{3}=0$ (the other cases are analogous). Our situation here is

$$
\begin{aligned}
& 0 \in a_{2} x_{2}+a_{3} x_{3} \\
& 0 \in b_{1} x_{1}+b_{3} x_{3}+b_{4} x_{4} \\
& 0 \in c_{1} x_{1}+c_{2} x_{2}+c_{4} x_{4}
\end{aligned}
$$

with all these coefficients different from zero. Then, multiplying the first, second and third equation by $a_{2}^{-1}, b_{1}^{-1}$ and $c_{1}^{-1}$ respectively, we get

$$
\begin{aligned}
& 0 \in x_{2}+a_{2}^{-1} a_{3} x_{3} \\
& 0 \in x_{1}+b_{1}^{-1} b_{3} x_{3}+b_{1}^{-1} b_{4} x_{4} \\
& 0 \in x_{1}+c_{1}^{-1} c_{2} x_{2}+c_{1}^{-1} c_{4} x_{4}
\end{aligned}
$$

Then we can suppose without loss of generality that $a_{2}=b_{1}=c_{1}=1$, and our situation is now

$$
\begin{aligned}
& 0 \in x_{2}+a_{3} x_{3} \\
& 0 \in x_{1}+b_{3} x_{3}+b_{4} x_{4} \\
& 0 \in x_{1}+c_{2} x_{2}+c_{4} x_{4}
\end{aligned}
$$

Pick $d_{1}=1, d_{2}=-c_{2}^{-1}, d_{3}=a_{3}^{-1} c_{2}^{-1}$ and $d_{4}=-b_{4}^{-1} a_{3}^{-1} b_{3} c_{2}^{-1}$. We have (using the fact that $F$ is hyperbolic) that

$$
d_{2}+a_{3} d_{3}=-c_{2}^{-1}+a_{3}\left(a_{3}^{-1} c_{2}^{-1}\right)=c_{2}^{-1}-c_{2}^{-1}=F
$$

and

$$
\begin{aligned}
d_{1}+b_{3} d_{3}+b_{4} d_{4} & =1+b_{3}\left(a_{3}^{-1} c_{2}^{-1}\right)+b_{4}\left(-b_{4}^{-1} a_{3}^{-1} b_{3} c_{2}^{-1}\right) \\
& =1+a_{3}^{-1} b_{3} c_{2}^{-1}-a_{3}^{-1} b_{3} c_{2}^{-1}=F
\end{aligned}
$$

and finally

$$
\begin{aligned}
d_{1}+c_{2} d_{2}+c_{4} d_{4} & =1+c_{2}\left(-c_{2}^{-1}\right)+c_{4}\left(-b_{4}^{-1} a_{3}^{-1} b_{3} c_{2}^{-1}\right) \\
& =1-1-a_{3}^{-1} b_{3} b_{4}^{-1} c_{2}^{-1} c_{4}=F .
\end{aligned}
$$

Then,

$$
\begin{aligned}
& 0 \in F=d_{2}+a_{3} d_{3} \\
& 0 \in F=d_{1}+b_{3} d_{3}+b_{4} d_{4} \\
& 0 \in F=d_{1}+c_{2} d_{2}+c_{4} d_{4},
\end{aligned}
$$

which prove that $\left(d_{1}, d_{2}, d_{3}, d_{4}\right)$ is a solution for this case.

Then fourth subcase is the one where $0 \notin\left\{a_{j}, b_{j}, c_{j}\right\}$ for some $j$. Suppose without loss of generality that $0 \notin\left\{a_{1}, b_{1}, c_{1}\right\}$. Since $F$ is a hyperfield, to find a non trivial weak solution of 3.8 is equivalent to find a non trivial weak solution of

$$
\begin{aligned}
& 0 \in x_{1}+a_{1}^{-1} a_{2} x_{2}+a_{1}^{-1} a_{3} x_{3}+a_{1}^{-1} a_{4} x_{4} \\
& 0 \in x_{1}+b_{1}^{-1} b_{2} x_{2}+b_{1}^{-1} b_{3} x_{3}+b_{1}^{-1} b_{4} x_{4} \\
& 0 \in x_{1}+c_{1}^{-1} c_{2} x_{2}+c_{1}^{-1} c_{3} x_{3}+c_{1}^{-1} c_{4} x_{4}
\end{aligned}
$$

Then we can suppose without loss of generality that $a_{1}=b_{1}=c_{1}=1$ and only deal with the new system

$$
\begin{align*}
& 0 \in x_{1}+a_{2} x_{2}+a_{3} x_{3}+a_{4} x_{4} \\
& 0 \in x_{1}+b_{2} x_{2}+b_{3} x_{3}+b_{4} x_{4} \\
& 0 \in x_{1}+c_{2} x_{2}+c_{3} x_{3}+c_{4} x_{4} \tag{7}
\end{align*}
$$

Note that, even in this reduced system we can suppose that for $j=2,3,4$ we have at most one element in $\left\{a_{j}, b_{j}, c_{j}\right\}$ equal to zero (because if one of the sets $\left\{a_{2}, b_{2}, c_{2}\right\},\left\{a_{3}, b_{3}, c_{3}\right\},\left\{a_{4}, b_{4}, c_{4}\right\}$ has two elements equal to zero, we are in the subcase two of the case III!). Then suppose without loss of generality that $a_{2} \neq 0, b_{3} \neq 0$ and $c_{4} \neq 0$.

Here we use again the fact that $F$ is hyperbolic: more specifically, we use that every hyperbolic hyperfield is rooted, in the sense that $\{a, b\} \subseteq a+b$ for all $a, b \in F^{*}$. Choose $d_{2}=-a_{2}^{-1}, d_{3}=-b_{3}^{-1}$ and $d_{4}=-c_{4}^{-1}$. Since $-1=a_{2} d_{2}=b_{3} d_{3}=b_{4} d_{4}$ we have

$$
\begin{aligned}
& -1 \in a_{2} d_{2}+a_{3} d_{3}+a_{4} d_{4} \\
& -1 \in b_{2} d_{2}+b_{3} d_{3}+b_{4} d_{4} \\
& -1 \in c_{2} d_{2}+c_{3} d_{3}+c_{4} d_{4}
\end{aligned}
$$

Picking now $d_{1}=1$, we have

$$
\begin{aligned}
& 1-1 \subseteq d_{1}+a_{2} d_{2}+a_{3} d_{3}+a_{4} d_{4} \\
& 1-1 \subseteq d_{1}+b_{2} d_{2}+b_{3} d_{3}+b_{4} d_{4} \\
& 1-1 \subseteq d_{1}+c_{2} d_{2}+c_{3} d_{3}+c_{4} d_{4}
\end{aligned}
$$

and then, $\left(d_{1}, d_{2}, d_{3}, d_{4}\right)$ is a non-trivial solution of the systems 3.9, which complete the proof for Case III. We further refer to this tactic to find $\left(d_{1}, d_{2}, d_{3}, d_{4}\right)$ as "the Case III method"l.

Case IV - the general case $m>n$ (and $n \geq 3$ ). Just proceed by induction on $m$. The base cases are Case I and II and the induction step is an argument similar to the the Case III method (3.8).

As application of these fragment of linear algebra for superfields, we get the following Theorem, which is a consequence of combining Theorem 3.6.12, Proposition 3.8.2, Theorem 3.8.18 and Theorem 3.8.23.

Theorem 3.8.24. Let $F$ be a linearly closed superfield and $p \in F[X]$ be an irreducible polynomial with $\operatorname{deg} p=n+1$. Then $F(p)$ is a full $F$-vector space and $\operatorname{dim}(F(p))=n+1$.

```
teolnc
```

Theorem 3.8.25. Let $F$ is a linearly closed superfield and $p \in F[X]$ be an irreducible polynomial with $\operatorname{deg} p=n+1$. Then $F(p)$ is also linearly closed.

Proof. Remember that $F(p)$ is generated by $\left\{1, \gamma, \ldots, \gamma^{n}\right\}$ with $\gamma=\bar{X}, n \in \mathbb{N}$. Also, we can consider $n$ as the minimal integer such that there exist $a_{0}, \ldots, a_{n+1}$ with

$$
0 \in d_{0}+d_{1} \gamma+\ldots+d_{n+1} \gamma^{n+1}
$$

Let $A \in M_{m \times q}(F(p))$, saying, $A=\left(\alpha_{i j}\right)$. We can write each $\alpha_{i j}$ as

$$
\alpha_{i j}=a_{0 i j}+a_{1 i j} \gamma+\ldots+a_{n i j} \gamma^{n}
$$

for suitable $a_{k i j} \in F$. Then a system $A x=0$ over $F(p)$ can be split into $n+1$ systems $A_{k} x=0$ over $F$, where $A_{k}=\left(a_{k i j}\right)$ for each $k=0,1, \ldots, n$ (in fact, $A x=0$ means $A_{0} x+\gamma A_{1} x+\gamma^{2} A_{2} x+\ldots+$ $\gamma^{n} A_{n} x=0$ ). Since $F$ is linearly closed, each $A_{k} x=0$ has at least a non-trivial solution, providing a non-trivial solution for $A x=0$.

As we can see, we had a lot of effort in order to prove Theorem 3.8.23. In this sense, we propose the following questions.

As we can see, we had a lot of effort in order to prove Theorem 3.8.23. In this sense, we propose the following questions.

## Question 3.8.26.

1. Is every hyperfield $F$ a linearly closed superfield?
2. Is every full superfield $F$ a linearly closed superfield?
3. What are the necessary conditions for a superfield $F$ be a linearly closed one?

In the context of algebraic and abstract theories of quadratic forms, there are at least two interesting Corollaries obtained applying Theorem 3.8.23 to the hyperfield $M(F):=F / m\left(F^{2} \backslash\{0\}\right)$ where $F$ is a field of characteristic not 2 , or more generally, $M(G)$ for a formally real special group $G$ (for a deeper understanding of $M(F)$ and $M(G)$, the reader can consult [47, [24], [23], [12], [17] or [45]).
cor-01
Corollary 3.8.27 (Isotropy Interpolation). Let $K=M(F):=F / m\left(F^{2} \backslash\{0\}\right)$ for a field $F$ (of characteristic not 2) or $K=M(G)$ for a formally real special group $G$. Consider a matrix $A \in$ $M_{n \times m}(K)$, saying $A=\left(a_{i j}\right)$. If $m>n$, there exists $d_{1}, \ldots, d_{n} \in F$, not all zero, such that all the forms $\left\{\varphi_{1}, \ldots, \varphi_{n}\right\}$ with

$$
\varphi_{i}:=\left\langle a_{i 1} d_{1}, a_{i 2} d_{2}, \ldots, a_{i m} d_{m}\right\rangle
$$

are isotropic.
cor-02
Corollary 3.8.28 (Hyperbolic Interpolation). Let $K=M(F) / m\left(M(F)^{2} \backslash\{0\}\right)$ where $M(F):=$ $F / m\left(F^{2} \backslash\{0\}\right)$ for a field $F$ (of characteristic not 2) or $K=M(G)$ for a formally real reduced special group $G$. Consider a matrix $A \in M_{n \times m}(K)$, saying $A=\left(a_{i j}\right)$. If $m>n$ is even, there exists $d_{1}, \ldots, d_{n} \in F$, not all zero, such that all the forms $\left\{\varphi_{1}, \ldots, \varphi_{n}\right\}$ with

$$
\varphi_{i}:=\left\langle a_{i 1} d_{1}, a_{i 2} d_{2}, \ldots, a_{i m} d_{m}\right\rangle
$$

are hyperbolic.
Also in the context of abstract theories of quadratic forms, Isotropic and Hyperbolic Interpolations 3.8 .27 and 3.8 .28 ) suggests interesting questions:

## Question 3.8.29.

1. In Corollaries 3.8 .27 and 3.8 .28 , are we able to get $d_{1}, \ldots, d_{n} \in F$, not all zero, such that all the Pfister forms $\left\{\varphi_{1}, \ldots, \varphi_{n}\right\}$ with

$$
\varphi_{i}:=\left\langle\left\langle a_{i 1} d_{1}, a_{i 2} d_{2}, \ldots, a_{i m} d_{m}\right\rangle\right\rangle
$$

are hyperbolic?
2. Are we able to get Corollary 3.8.28 for general fields or general special-groups (not necessarily reduced)?

As application of these fragment of linear algebra for superfields, we get the following Theorem, which is a consequence of combining Theorem 3.6.12, Proposition 3.8.2, Theorem 3.8.18 and Theorem 3.8.23,

Theorem 3.8.30. Let $F$ be a linearly closed superfield and $p \in F[X]$ be an irreducible polynomial with $\operatorname{deg} p=n+1$. Then $F(p)$ is a full $F$-vector space, with a rigid basis $\left\{1, \gamma, \ldots, \gamma^{n}\right\} \quad(\gamma:=[X] \in$ $F[X] /\langle p(X)\rangle)$ and $\operatorname{dim}(F(p))=n+1$.

### 3.9 A quantifier elimination procedure

We also have a quantifier elimination procedure for any infinite algebraically closed associative superfield. This is a variation of Theorem 9.2.1 in [36] and a generalization of the results in [19].

Throughout this Section, all superfields will be considered associative.
Lemma 3.9.1 (Lemma 1.27 of (19]). Let $A$ be a superring.
$i$ - For all $n \in \mathbb{N}$ and all $a_{0}, \ldots, a_{n-1} \in A$, the sum $a_{0}+\ldots+a_{n-1}$ and product $a_{0} \cdot \ldots \cdot a_{n-1}$ does not depends on the order of the entries.
ii - For every term $t\left(y_{1}, \ldots, y_{n}\right)$ on the 2-ring language, exists variables $x_{i j}$ such that $A$ satisfies the formula

$$
t\left(y_{1}, \ldots, y_{n}\right) \sqsubseteq \sum_{i<p} \prod_{j<m_{i}} x_{i j} .
$$

Moreover, if $A$ is a full 2-ring, it satisfies the formula

$$
t\left(y_{1}, \ldots, y_{n}\right)={ }_{s} \sum_{i<p} \prod_{j<m_{i}} x_{i j} .
$$

Lemma 3.9.2 (Lemma 3.2 of [19]). Let A be a superring, $t_{1}(\bar{x}), t_{2}(\bar{x})$ be terms on the full superring language and let $v=\bar{a}: \bar{x} \rightarrow A$

$$
i-t_{1}^{A}(\bar{a}) \subseteq t_{2}^{A}(\bar{a}) \text { iff } 0 \in\left(t_{2}-t_{1}\right)^{A}(\bar{a}) .
$$

ii - Given any atomic formula, $t_{1}(\bar{x}) \sqsubseteq t_{2}(\bar{x})$, there is a polynomial term $p(\bar{x}) \in R[\bar{x}]$ such that

$$
A \models_{v}\left(t_{1}(\bar{x}) \sqsubseteq t_{2}(\bar{x})\right) \leftrightarrow(0 \sqsubseteq p(\bar{x})) .
$$

Let $\mathcal{L}$ be the language of superrings. For each superring $R$, let $\mathcal{L}(R)$ be the language extending $\mathcal{L}$ by adding all elements of $R$ as strict constant symbols. Let $\Gamma^{\prime}$ be the superring axioms. Let extend $\Gamma^{\prime}$ by (in)equalities and relations of the form

$$
a_{0} \neq b_{0} ; c_{1}=a_{1} . b_{1} ; c_{2} \in a_{2}+b_{2} ; a_{i}, b_{i}, c_{i} \in R
$$

that are true in $R$ ("the diagram of $R$ "). Denote the set of formulas obtained by $\Gamma^{\prime}(R)$. A model of $\Gamma^{\prime}(R)$ is a superring that contains a subset $\bar{R}=\{\bar{a}: a \in R\}$ and $\bar{R}$ is an isomomorphic copy of $R$ inside this model. If $R=K$ is a superfield and $\Gamma$ is the superfield axioms, then a model of $\Gamma(K)$ is a superfield that contains a subset $\bar{K}=\{\bar{a}: a \in K\}$ and $\bar{K}$ is a superfield isomorphic to $K$. Then a model of $\Gamma(K)$ is (up to a isomorphism) a superfield containing $K$. Now, we extend $\Gamma(K)$ to a new set of axioms $\tilde{\Gamma}(K)$ adding axioms to obtain an algebraic closure superfield

$$
\begin{equation*}
\forall z_{0} \ldots \forall z_{n} \exists x\left[0 \in z_{0}+z_{1} x+\ldots+z_{n-1} x^{n-1}+x^{n}\right], n \geq 1 . \tag{AC}
\end{equation*}
$$

We add also the family of axioms $\exists z_{0} \ldots \exists z_{n-1} \underset{i<j<n}{\bigvee}\left[z_{i} \neq z_{j}\right], n \geq 2$.
A model $F$ of $\Gamma(K)$ is also a model of $\tilde{\Gamma}(K)$ iff $F$ is infinite and algebraically closed. Our aim is to describe a quantifier elimination procedure for $\tilde{\Gamma}(F)$. By the reduction Lemma 3.9.2, $F$ regards every atomic formula as equivalent modulo $\Gamma(K)$ to a polynomial "equation" $0 \in f\left(X_{1}, \ldots, X_{n}\right)$. Since $K[\bar{X}]$ is a superdomain, a conjunction of inequations $\bigwedge_{i=1}^{m}\left[0 \neq g_{i}(\bar{X})\right]$ is equivalent to the "inequation" $0 \notin g_{1}(\bar{X}) \ldots g_{n}(\bar{X})$. Then, to obtain a quantifier elimination for $\tilde{\Gamma}(K)$ is sufficient eliminate $Y$ from the formula

$$
\begin{equation*}
\exists Y\left[0 \in f_{1}(\bar{X}, Y) \wedge \ldots \wedge 0 \in f_{m}(\bar{X}, Y) \wedge 0 \notin g(\bar{X}, Y)\right] \tag{3.910}
\end{equation*}
$$

with $f_{1}, \ldots, f_{m}, g \in R\left[X_{1}, \ldots, X_{m}, Y\right]$.
Theorem 3.9.3 (Quantifier Elimination Procedure, Adapted from Theorem 3.3 of [19]). Let $K$ be an infinite superfield and $\varphi\left(X_{1}, \ldots, X_{n}, Y\right)$ the formula in 3.10. Then $\varphi\left(X_{1}, \ldots, X_{n}, Y\right)$ is equivalent modulo $\tilde{\Gamma}(R)$ to a Boolean combination of atomic formulas $\psi\left(X_{1}, \ldots, X_{r}\right), r \geq n$.

Proof. The proof consists in three parts:
A - Reduction to the case that only one of $f_{1}, \ldots, f_{m}$ involves $Y$. Move each conjunction that appears in (3.10) and that does not involve $Y$ to the left of $\exists Y$ according to the rule " $\exists Y[\varphi \wedge \psi] \equiv$ $\varphi \wedge \exists Y[\psi]$ if $Y$ does not appear in $\varphi^{\prime \prime}$. Thus we assume $\operatorname{deg}_{Y}\left(f_{i}(\bar{X}, Y)\right) \geq 1, i=1, \ldots, m$ and $m \geq 2$. We now perform an induction on $\sum \operatorname{deg}_{Y}\left(f_{i}(\bar{X}, Y)\right)$ : Let $p(\bar{X}, Y)$ and $q(\bar{X}, Y)$ be multipolynomials with coefficients in $R$ such that $0 \leq \operatorname{deg}_{Y} p(\bar{X}, Y) \leq \operatorname{deg}_{Y} q(\bar{X}, Y)=d$. Write $p(\bar{X}, Y)$ in the form

$$
\begin{equation*}
p(\bar{X}, Y)=a_{k}(\bar{X}) Y^{k}+a_{k-1}(\bar{X}) Y^{k-1}+\ldots+a_{0}(\bar{X}) \tag{9.9}
\end{equation*}
$$

with $a_{j} \in R[\bar{X}]$. For each $j$ with $0 \leq j \leq k$ let

$$
p_{j}(\bar{X}, Y)=a_{j}(\bar{X}) Y^{j}+a_{j-1}(\bar{X}) Y^{j-1}+\ldots+a_{0}(\bar{X})
$$

If $0 \notin a_{j}(\bar{X})$, division of $q(\bar{X}, Y)$ by $p_{j}(\bar{X}, Y)$ produces $q_{j}(\bar{X}, Y)$ and $r_{j}(\bar{X}, Y)$ in $R[\bar{X}, Y]$ for which

$$
\begin{equation*}
a_{j}(\bar{X})^{d} q(\bar{X}, Y) \subseteq q_{j}(\bar{X}, Y) p_{j}(\bar{X}, Y)+r_{j}(\bar{X}, Y), \tag{3.92}
\end{equation*}
$$

and $\operatorname{deg}_{Y}\left(r_{j}\right)<\operatorname{deg}_{Y}\left(p_{j}\right) \leq d$. Let $F$ be a model of $\Gamma(K)$. If $x_{1}, \ldots, x_{n}, y$ are elements of $F$ such that $0 \in a_{l}(\bar{x})$ for $l=j+1, \ldots, k$ and $0 \notin a_{j}(\bar{x})$, then $[0 \in p(\bar{x}, y) \wedge 0 \in q(\bar{x}, y)]$ is equivalent in $F$ to $\left[0 \in p_{j}(\bar{x}, y) \wedge 0 \in r_{j}(\bar{x}, y)\right]$. Therefore, the formula $[0 \in p(\bar{X}, Y) \wedge 0 \in q(\bar{X}, Y)]$ is equivalent modulo $\Gamma(K)$ to the formula

$$
\begin{align*}
& \left(\bigvee_{j=0}^{k}\left[0 \in a_{k}(\bar{X}) \wedge \ldots \wedge 0 \in a_{j+1}(\bar{X}) \wedge 0 \notin a_{j}(\bar{X}) \wedge 0 \in p_{j}(\bar{X}, Y) \wedge 0 \in r_{j}(\bar{X}, Y)\right]\right) \\
& \vee\left[0 \in a_{k}(\bar{X}) \wedge \ldots \wedge 0 \in a_{0}(\bar{X}) \wedge 0 \in q(\bar{X}, Y)\right] . \tag{3.934}
\end{align*}
$$

Apply the outcome of (3.13) to $f_{1}(\bar{X}, Y)$ and $f_{m}(\bar{X}, Y)$ (of 3.10 ). With the rule " $\exists Y[\varphi \vee \psi] \equiv$ $\exists Y \varphi \vee \exists Y \psi$ " we have replaced (3.10) by disjunction of statements of form (3.10) in each which the sum corresponding to $\sum \operatorname{deg}_{Y}\left(f_{i}(X, Y)\right)$ is smaller. Using the induction assumption, we conclude that $m$ may be taken to be at most 1 .

B - Reduction to the case that $m=0$. Continue the notation of part $A$ which left us at the point of considering how to eliminate $Y$ from $p(\bar{X}, Y)$ in

$$
\begin{equation*}
\exists Y[0 \in p(\bar{X}, Y) \wedge 0 \notin g(\bar{X}, Y)] . \tag{2}
\end{equation*}
$$

Consider a model $F$ of $\tilde{\Gamma}(K)$ and elements $x_{1}, \ldots, x_{n} \in F$. If $0 \notin p(\bar{x}, Y)$ then (since $F$ is algebraically closed) the statement " $F \models \exists Y[0 \in p(\bar{x}, Y) \wedge 0 \notin g(\bar{x}, Y)]$ " is equivalent to the statement $" p(\bar{x}, Y)$ does not divide $g(\bar{x}, Y)^{k}$ in $F[X] "$. Therefore, with $q(\bar{X}, Y)=g(\bar{X}, Y)^{k}$ and in the notation of (3.11) and 3.12, formula (3.14) is equivalent modulo $\tilde{\Gamma}(K)$ to the formula

$$
\begin{array}{r}
\left(\bigvee_{j=0}^{k}\left[0 \in a_{k}(\bar{X}) \wedge \ldots \wedge 0 \in a_{j+1}(\bar{X}) \wedge 0 \notin a_{j}(\bar{X}) \wedge \exists Y\left[\in r_{j}(\bar{X}, Y)\right]\right]\right) \\
\vee\left[0 \in a_{k}(\bar{X}) \wedge \ldots \wedge 0 \in a_{0}(\bar{X}) \wedge \exists Y[0 \in g(\bar{X}, Y)]\right]
\end{array}
$$

a disjunction of statements of form (3.10) with $m=0$.

C - Completion of the proof. By part B we are in the point of removing $Y$ from a statement of the form $\exists Y\left[0 \notin a_{l}(\bar{X}) Y^{l}+a_{l-1}(\bar{X}) Y^{l-1}+\ldots+a_{0}(\bar{X})\right]$. Since models of $\tilde{\Gamma}(K)$ are infinite superfields, this formula is equivalent modulo $\tilde{\Gamma}(K)$ to $0 \notin a_{l}(\bar{X}) \vee \ldots \vee 0 \notin a_{0}(\bar{X})$, completing the quantifier elimination procedure.

## Chapter 4

# K-theories: the rise of (universal) Inductive Graded Rings 

Concerning Abstract Theories of Quadratic forms (in particular special groups and real semigroups), the references [28], [32] and [33] are central. The theory of special groups deals simultaneously reduced and non-reduced theories but focuses on rings with an "expressive amount" of invertible coefficients to quadratic forms and the theory of real semigroups consider general coefficients of a ring, but only addresses the reduced case. Both are first-order theory, thus they allow the use of model theoretic methods.
M. Marshall in [47] introduced an approach to (reduced) theory of quadratic forms through the concept of multiring ${ }^{1}$ : this seems more intuitive for an algebraist since it encompasses (generalizes, in fact) some techniques of ordinary Commutative Algebra. Moreover, the multirings encode copies of special groups and real semigroups (see [24]) and still allows the use of model-theoretic tools, since multirings (hyperrings) endowed with convenient notion of morphisms constitutes a category that is isomorphic to a category of appropriate first-order structures.

In the recent work [17]: (i) we have considered interesting pairs $(A, T)$ where $A$ is a multiring and $T \subseteq A$ is a certain multiplicative subset in such a way to obtain models of abstract theories of quadratic forms (special groups and real semigroups) via natural quotients - Marshall's quotient construction and (ii) we have used this new setting to motivate a "non reduced" expansion of the theory of real semigroups to deal the formally real case, isolating axioms over pairs involving multirings and a multiplicative subset with some properties.

The uses of K-theoretic (and Boolean) methods in abstract theories of quadratic forms has been proved a very successful method, see for instance, these two papers of Dickmann and Miraglia: [27] where they give an affirmative answer to Marshall's Conjecture, and [29], where they give an affirmative answer to Lam's Conjecture.

These two central papers makes us take a deeper look at the theory of Special Groups by itself. This is not mere exercise in abstraction: from Marshall's and Lam's Conjecture many questions arise in the abstract and concrete context of quadratic forms. Even in the algebraic theory of quadratic forms, there are simple (and unsolved questions), some of them solved just in the last decade, as showed by 40.

With these two paragraphs in mind, the purpose of this Chapter is to prepare the land for further generalizations (with applications) of the "Milnor's triangle K-theory - quadratic forms

[^12]- Galois cohomology". The main results are Theorems 4.5.6 and its Corollaries, which provides interchanging formulas between the three K-theories considered here.


### 4.1 Milnor's K-theory

For further references, in this section we get some definitions and results about Milnor's Ktheory, as developed in [52. Before deal with Milnor's K-theory, lets make a brief summary on graded rings.

Definition 4.1.1 (Graded Ring [9]). Let ( $G, \cdot$ ) be a monoid. A ring $A$ is said to be $G$-graded if its additive group $(A,+)$ admits a decomposition in direct sum of abelian groups

$$
A=\bigoplus_{g \in G} A_{g},
$$

satisfying $A_{g} \cdot A_{h} \subseteq A_{g \cdot h}$ for all $g, h \in G$, or in other words,

$$
a_{g} \in A_{g}, a_{h} \in A_{h} \Rightarrow a_{g} a_{h} \in A_{g \cdot h}(g, h \in G) .
$$

The elements $a_{g} \in A_{g} \subseteq A$ are called homogeneous of degree $g$. Then, every element of $A$ can be written uniquely as a sum $a=\sum_{g \in G} a_{g}$ of homogeneous elements $a_{g} \in A_{g}$. We call $a_{g}$ of homogeneous component of degree $g$ of the element $a$.

A morphism of $G$-graded rings is a ring homomorphism $\varphi: A \rightarrow B$ that respect the graduation, i.e, such that for all $g \in G, \varphi\left(A_{g}\right) \subseteq B_{g}$. The category of $G$-graded rings and its morphisms will be denoted by Grad $_{G}$.

The most important cases are those when $G=\mathbb{Z}$ or $G=\mathbb{N}$. Since a $\mathbb{N}$-graded ring can be seen as a $\mathbb{Z}$-graded ring with components of negative degree equal to zero, unless we mention, we call the $\mathbb{Z}$-graded rings just by graded rings and we will denote grad := $\operatorname{grad}_{\mathbb{Z}}$.

Example 4.1.2. Let $A$ be a ring. The "canonical" example of graded ring is $A\left[x_{1}, \ldots, x_{n}\right]$, that admits a graduation

$$
A\left[x_{1}, \ldots, x_{n}\right]=\bigoplus_{d \geq 0} A\left[x_{1}, \ldots, x_{n}\right]_{d}
$$

when $A\left[x_{1}, \ldots, x_{n}\right]_{d}$ is the free $A$-module of rank $\binom{n+d-1}{d}$ with basis given by the monomials $x_{1}^{e_{1}} x_{2}^{e_{2}} \ldots x_{n}^{e_{n}}$ of degree $d=e_{1}+\ldots+e_{n}$.

Definition 4.1.3 (Homogeneous Ideal). Let $A$ be a $G$-graded ring. $A n$ ideal $I \subseteq A$ is said to be an homogeneous ideal if $(I,+)$ admits a decomposition

$$
I=\bigoplus_{g \in G}\left(I \cap A_{g}\right)
$$

Lemma 4.1.4. Let $(G,+)$ be an abelian group, $A=\bigoplus_{g \in G} A_{g}$ a $G$-graded ring and $I \subseteq A$ an ideal. For each element $a \in A$, denotes by $a_{g} \in A_{g}$ its homogeneous component of degree $g$. Are equivalent:

$$
i \text { - } I \text { is an homogeneous ideal; }
$$

ii - for all $a \in A$,

$$
a \in I \Leftrightarrow a_{g} \in I \text { for all } g \in G ;
$$

iii - I is generated by homogeneous elements (possibly of different degrees).
Proof. (i) $\Leftrightarrow$ (ii) and (ii) $\Rightarrow$ (iii) are direct consequences of the definitions involved in. For (iii) $\Rightarrow$ (ii), suppose that $I$ is generated by homogeneous elements $a_{i}(i$ in some set of index $\Lambda)$ and let $a \in I$. Then we can write $a=b_{1} a_{i_{1}}+\ldots+b_{n} a_{i_{n}}$ with $b_{i} \in A=\bigoplus_{g \in G} A_{g}$. Expanding each $b_{i}=\sum_{g \in G} b_{i g}$ as sum of its homogeneous components in $b_{i g} \in A_{g}$, the degree $g$ term of $a$ is

$$
a_{g}=b_{1\left(g-\operatorname{deg}\left(a_{i_{1}}\right)\right.} a_{i_{1}}+. .+b_{n\left(g-\operatorname{deg}\left(a_{i_{n}}\right)\right.} a_{i_{n}} \in I
$$

Lemma 4.1.5 (Generating Homogeneous Ideals). Let $A=\bigoplus_{n \geq 0} A_{n}$ be a graded ring. Let $a \in A_{0}$. Consider the following recursive construction:

$$
\begin{gathered}
\left.I_{0}:=\langle a\rangle \subseteq A_{0} \text { (the ideal generated by a on } A_{0}\right) \\
I_{n}:=\left\langle x \cdot y: x \in I_{p}, y \in I_{q} \text { with } p+q=n\right\rangle \subseteq A_{n} .
\end{gathered}
$$

Then $I=\left(I_{n}\right)_{n \geq 0}$ is an homogeneous ideal of $A$, called the homogeneous ideal generated by $a$.

So lets present the basic definitions and properties of Milnor's K-theory (as described in [52]). milkt
Definition 4.1.6 (The Milnor's K-theory of a Field [52]). For a field F (of characteristic not 2), $K_{*} F$ is the graded ring

$$
K_{*} F=\left(K_{0} F, K_{1} F, K_{2} F, \ldots\right)
$$

defined by the following rules: $K_{0} F:=\mathbb{Z} . K_{1} F$ is the multiplicative group $\dot{F}$ written additively. With this purpose, we fix the canonical"logarithm" isomorphism

$$
l: \dot{F} \rightarrow K_{1} F
$$

where $l(a b)=l(a)+l(b)$. Then $K_{n} F$ is defined to be the quotient of the tensor algebra

$$
K_{1} F \otimes K_{1} F \otimes \ldots \otimes K_{1} F(n \text { times })
$$

by the (homogeneous) ideal generated by all $l(a) \otimes l(1-a)$, with $a \neq 0,1$. We also have the reduced $K$-theory graded ring $k_{*} F=\left(k_{0} F, k_{1} F, \ldots, k_{n} F, \ldots\right)$, which is defined by the rule $k_{n} F:=K_{n} F / 2 K_{n} F$ for all $n \geq 0$.

With these definitions, the K-theory structure gives us the following three basic Lemmas:
Lemma 4.1.7 (1.1 [52]). For every $\xi \in K_{m} F$ and $\eta \in K_{n} F$, the identity

$$
\eta \xi=(-1)^{m n} \xi \eta
$$

is valid in $K_{n+m} F$.
Lemma 4.1.8 (1.2 [52]). The identity $l(a) \otimes l(a)=l(a) \otimes l(-1)$ is valid for every $l(a) \in K_{1} F$.
Lemma 4.1.9 (1.3 [52]). If the sum $a_{1}+\ldots+a_{n}$ of non-zero field elements is equal to either 0 or 1 , then $l\left(a_{1}\right) \otimes \ldots \otimes l\left(a_{n}\right)=0$.

Theorem 4.1.10 (Theorem 4.1 of [52]). there is only one morphism

$$
s_{n}: k_{n} F \rightarrow I^{n} F / I^{n+1} F
$$

which carries each product $l\left(a_{1}\right) \ldots l\left(a_{n}\right)$ in $K_{n} F / 2 K_{n} F$ to the product $\left(\left\langle a_{1}\right\rangle-\langle 1\rangle\right) \ldots\left(\left\langle a_{n}\right\rangle-\langle 1\rangle\right)$ modulo $I^{n+1} F$.

These morphisms determines a surjective $s_{*}: k_{*} F \rightarrow W_{g}(F)$, where

$$
W_{g}(F)=\left(W F / I F, I F / I^{2} F, \ldots, I^{n} F / I^{n+1} F, \ldots\right)
$$

For a field $F$, let $F_{s}$ be the a separable closure of $F$ and $G_{F}=\operatorname{Gal}\left(F_{s}\right)$. Then, the exact sequence

$$
1 \longrightarrow\{ \pm 1\} \longrightarrow \dot{F}_{s} \xrightarrow{2} \dot{F}_{s} \longrightarrow 1
$$

is taken to the following exact sequence

$$
H^{0}\left(G_{F}, \dot{F}_{s}\right) \xrightarrow{2} H^{0}\left(G_{F}, \dot{F}_{s}\right) \xrightarrow{\delta} H^{1}\left(G_{F},\{ \pm 1\}\right) \longrightarrow H^{1}\left(G_{F}, \dot{F}_{s}\right)
$$

of cohomology groups. Identifying the two first groups with $\dot{F}$, and $\{ \pm 1\}$ with $\mathbb{Z} / 2 \mathbb{Z}$ and applying Hilbert's 90 , we have

$$
\dot{F} \xrightarrow{2} \dot{F} \xrightarrow{\delta} H^{1}\left(G_{F}, \mathbb{Z} / 2 \mathbb{Z}\right) \longrightarrow 0
$$

The quotient $\dot{F} / \dot{F}^{2}$ is identified with $k_{1} F$.
Theorem 4.1.11 (Lemma 6.1 of [52]). The isomorphism $l(a) \mapsto \delta(a)$ from $K_{1} F / 2 K_{1} F$ to $H^{1}\left(G_{F}, \mathbb{Z} / 2 \mathbb{Z}\right)$ admits a unique extension to a graded ring morphism

$$
h_{f}: k_{*} F \rightarrow H^{*}\left(G_{F}, \mathbb{Z} / 2 \mathbb{Z}\right) .
$$

Milnor's Conjecture consists to say that $s$ and $h$ are graded rings isomorphisms, which makes the fuctors $K_{*} F / 2 K_{*} F, W_{g}(F), H^{*}(G, \mathbb{Z} / 2 \mathbb{Z})$ isomorphic.

### 4.2 Dickmann-Miraglia K-theory for Special Groups

There are some generalizations of Milnor's K-theory. In the quadratic forms context, maybe the most significant one is the Dickmann-Miraglia K-theory of Special Groups. It is a main tool in the proof of Marshall's and Lam's Conjecture. In this section, we get some definitions and results from [28] and [30].
defn:ksg
Definition 4.2.1 (The Dickmann-Miraglia K-theory [30]). For each special group $G$ (written multiplicatively) we associate a graded ring

$$
k_{*} G=\left(k_{0} G, k_{1} G, \ldots, k_{n} G, \ldots\right)
$$

as follow: $k_{0} G:=\mathbb{F}_{2}$ and $k_{1} G:=G$ written additively. With this purpose, we fix the canonical "logarithm" isomorphism $\lambda: G \rightarrow k_{1} G, \lambda(a b)=\lambda(a)+\lambda(b)$. Observe that $\lambda(1)$ is the zero of $k_{1} G$ and $k_{1} G$ has exponent 2, i.e, $\lambda(a)=-\lambda(a)$ for all $a \in G$. In the sequel, we define $k_{*} G$ by the quotient of the $\mathbb{F}_{2}$-graded algebra

$$
\left(\mathbb{F}_{2}, k_{1} G, k_{1} G \otimes_{\mathbb{F}_{2}} k_{1} G, k_{1} G \otimes_{\mathbb{F}_{2}} k_{1} G \otimes_{\mathbb{F}_{2}} k_{1} G, \ldots\right)
$$

by the (graded) ideal generated by $\left\{\lambda(a) \otimes \lambda(a b), a \in D_{G}(1, b)\right\}$. In other words, for each $n \geq 2$,

$$
k_{n} G:=T^{n}\left(k_{1} G\right) / Q^{n}(G),
$$

where

$$
T^{n}\left(k_{1} G\right):=k_{1} G \otimes_{\mathbb{F}_{2}} k_{1} G \otimes_{\mathbb{F}_{2}} \ldots \otimes_{\mathbb{F}_{2}} k_{1} G
$$

and $Q^{n}(G)$ is the subgroup generated by all expressions of type $\lambda\left(a_{1}\right) \otimes \lambda\left(a_{2}\right) \otimes \ldots \otimes \lambda\left(a_{n}\right)$ such that for some $i$ with $1 \leq i<n$, there exist $b \in G$ such that $a_{i} \in D_{G}(1, b)$ and $a_{i}=a_{i+1} b$, which in symbols, means

$$
\begin{aligned}
Q^{n}(G): & :=\left\langle\left\{\lambda\left(a_{1}\right) \otimes \lambda\left(a_{2}\right) \otimes \ldots \otimes \lambda\left(a_{n}\right): \text { exists } 1 \leq i<n \text { and } b \in G\right.\right. \\
& \text { such that } \left.\left.a_{i}=a_{i+1} b \text { and } a_{i} \in D_{G}(1, b)\right\}\right\rangle .
\end{aligned}
$$

Since $\lambda(a)+\lambda(a)=0$ for all $a \in k_{1} G$, follow that $\eta+\eta=0$ for all $\eta \in k_{n} G$, so this is a group of exponent 2. Moreover, for all $a, b \in G$,

$$
\overline{\lambda(a) \otimes \lambda(a b)}=\overline{\lambda(a) \otimes[\lambda(a)+\lambda(b)]}=\overline{\lambda(a) \otimes \lambda(a)+\lambda(a) \otimes \lambda(b)}=\overline{\lambda(a) \otimes \lambda(a)}+\overline{\lambda(a) \otimes \lambda(b)},
$$

hence

$$
a \in D_{G}(1, b) \Rightarrow \overline{\lambda(a) \otimes \lambda(a b)}=\overline{0} \text { in } k_{2} G,
$$

or equivalently, $\overline{\lambda(a) \otimes \lambda(a)}=\overline{\lambda(a) \otimes \lambda(b)}$ in $k_{2} G$.
Before we proceed, lets make some abbreviations/simplifications to make the reading of this work more easy and comfortable. Firstly, whenever possible, we will omit the over line that indicates the equivalence classes. For example, the affirmation " $\overline{\lambda(a) \otimes \lambda(a)}=\overline{\lambda(a) \otimes \lambda(b)}$ in $k_{2} G^{\prime \prime}$ will be expressible in the simplified manner by " $\lambda(a) \otimes \lambda(a)=\lambda(a) \otimes \lambda(b)$ in $k_{2} G$ ".

Moreover, we will denote " $\lambda\left(a_{1}\right) \otimes \lambda\left(a_{2}\right) \otimes \ldots \otimes \lambda\left(a_{n}\right)$ " simply by " $\lambda\left(a_{1}\right) \lambda\left(a_{2}\right) \ldots \lambda\left(a_{n}\right)$ ". Sure, $k_{*}(G)$ is a graded ring, and in particular, a ring, so that we are able to multiply elements $\eta \in k_{n} G$, $\tau \in k_{m} G$. Whenever we want to do this, we will denote $\eta \cdot \tau$, in order to avoid confusion with the simplifications described above.

Finally, since we only take tensorial products with parameters in $\mathbb{F}_{2}$, we abbreviate " $A \otimes_{\mathbb{F}_{2}} B$ " simply by " $A \otimes B$ ". In this way, $T^{n}\left(k_{1} G\right)$ we will be denoted simply by

$$
T^{n}\left(k_{1} G\right)=k_{1} G \otimes k_{1} G \otimes \ldots \otimes k_{1} G .
$$

Next, we have a result that approximate Dickmann-Miraglia's K-theory with the Milnor's reduced K-theory:
2.1kt

Proposition 4.2.2 (2.1 [30]). Let $G$ be a special group, $x, y, a_{1}, \ldots, a_{n} \in G$ and $\sigma$ be a permutation on $n$ elements.
$a-$ In $k_{2} G, \lambda(a)^{2}=\lambda(a) \lambda(-1)$. Hence in $k_{m} G, \lambda(a)^{m}=\lambda(a) \lambda(-1)^{m-1}, m \geq 2$;
$b-$ In $k_{2} G, \lambda(a) \lambda(-a)=\lambda(a)^{2}=0$;
$c-\operatorname{In} k_{n} G, \lambda\left(a_{1}\right) \lambda\left(a_{2}\right) \ldots \lambda\left(a_{n}\right)=\lambda\left(a_{\sigma 1}\right) \lambda\left(a_{\sigma 2}\right) \ldots \lambda\left(a_{\sigma n}\right)$;
$d$ - For $n \geq 1$ and $\xi \in k_{n} G, \xi^{2}=\lambda(-1)^{n} \xi$;
$e$ - If $G$ is a reduced special group, then $x \in D_{G}(1, y)$ and $\lambda(y) \lambda\left(a_{1}\right) \ldots \lambda\left(a_{n}\right)=0$ implies

$$
\lambda(x) \lambda\left(a_{1}\right) \lambda\left(a_{2}\right) \ldots \lambda\left(a_{n}\right)=0 .
$$

An element $a \in G$ induces a graded morphism of degree $1, \omega^{a}=\left\{\omega_{n}^{a}\right\}_{n \geq 1}: k_{*} G \rightarrow k_{*} G$, where $\omega_{n}^{a}: k_{n} G \rightarrow k_{n} G$ is the multiplication by $\lambda(-a)$. When $a=-1$, we write

$$
\omega=\left\{\omega_{n}\right\}_{n \geq 1}=\left\{\omega_{n}^{-1}\right\}_{n \geq 1}=\omega^{-1} .
$$

The Lemma 4.2 .3 below generalizes Proposition 5.10 of [59. Firstly, lets establish some notation. For $n \geq 0$ denote $P(n):=\mathcal{P}(\{0, \ldots, n-1\}) \backslash\{\emptyset\}$ and for $0 \leq i \leq n-1$, denote $P(n, i):=\{X \in P(n): i \in X\}$. Now, let $G$ be a pre-special group and $\left\{a_{0}, \ldots, a_{n-1}\right\} \subseteq G$ $\mathbb{F}_{2}$-linearly independent. If $S \in P(n)$ we denote

$$
a_{S}:=a_{0}^{\varepsilon_{0}} \ldots a_{n-1}^{\varepsilon_{n-1}},
$$

where $\varepsilon_{0} \in\{0,1\}$ for all $i=0, . ., n-1$ and $\varepsilon_{i}=1$ if and only if $i \in S$. Remember that by the very definition of $k_{n}(G)$,

$$
k_{n}(G):=\left[k_{1}(G) \otimes k_{1}(G)\right] / M,
$$

where $M$ is the subgroup of $k_{1}(G) \otimes k_{1}(G)$ generated by

$$
\left\{\lambda(a) \lambda(b): a \in D_{G}(1, b), a, b \in G\right\}
$$

Lemma 4.2.3. Let $G$ be a pre-special group and $\left\{a_{0}, \ldots, a_{n-1}\right\} \subseteq G \mathbb{F}_{2}$-linearly independent. Are equivalent:
$i$ - There exists $\left\{b_{0}, \ldots, b_{n-1}\right\} \subseteq G$ such that

$$
\sum_{k<n} \lambda\left(a_{k}\right) \lambda\left(b_{k}\right)=0 \text { in } k_{2}(G) .
$$

ii - There exists subsets $\left\{c_{0}, \ldots, c_{m-1}\right\},\left\{d_{0}, \ldots, d_{n-1}\right\}$ of $G$ with $m \geq n$ such that
(a) $\left\{c_{0}, \ldots, c_{m-1}\right\}$ is linearly independent and $c_{i}=a_{i}$ for all $i<n$;
(b) $d_{i}=b_{i}$ for all $i<n$ and $d_{i}=1$ for $i=n, \ldots, m-1$.
(c) For all $x \in C:=\left[c_{0}, \ldots, c_{m-1}\right]$, there is some $r_{x} \in D_{G}(1, x)$ such that for each $i<m$

$$
d_{i}=\prod_{x \in C_{i}} r_{x}
$$

where

$$
C_{i}=\left\{\prod_{k<m} c_{k}^{\varepsilon_{k}}: \varepsilon_{k} \in\{0,1\} \text { and } \varepsilon_{i}=1\right\} .
$$

In other words, $C_{i}$ is "counting" all products $c_{0}^{r_{0}} \ldots c_{i}^{1} \ldots c_{m-1}^{r_{m-1}}$. Since for all $x \in C:=\left[c_{0}, \ldots, c_{m-1}\right]$ there exist $S \in P(m)$ such that

$$
x=\prod_{i \in S} c_{i}:=c_{S} .
$$

Denoting $r_{x}$ by $r_{S}$ we can rewrite

$$
d_{i}=\prod_{x \in C_{i}} r_{x}=\prod_{S \in P(m)} r_{S} .
$$

Proof of Lemm\& 4.2.3. (i) $\Rightarrow$ (ii). Let

$$
\sum_{k<n} \lambda\left(a_{k}\right) \lambda\left(b_{k}\right)=0 \text { in } k_{2}(G) .
$$

Then there exist $u_{0}, \ldots, u_{p-1}, v_{0}, \ldots, v_{p-1} \in G$ such that $v_{i} \in D_{G}\left(1, u_{i}\right)$ for $i=0, \ldots, p-1$ and

$$
\sum_{k<n} \lambda\left(a_{k}\right) \lambda\left(b_{k}\right)=\sum_{k<p} \lambda\left(u_{k}\right) \lambda\left(v_{k}\right) \text { in } k_{1}(G) \otimes k_{1}(G) .
$$

Enlarge the set $\left\{a_{0}, \ldots, a_{n-1}\right\}$ to a base $\left\{c_{0}, \ldots, c_{m-1}\right\}$ of $\left[\left\{a_{0}, \ldots, a_{n-1}, u_{0}, \ldots, u_{p-1}\right\}\right]$, with $c_{i}=a_{i}$ for all $i<n$. For all $x \in C:=\left[c_{0}, \ldots, c_{m-1}\right]$ there exist an unique $S \in P(m)$ such that

$$
x=\prod_{i \in S} c_{i}:=c_{S} .
$$

Moreover, since $\left\{c_{0}, \ldots, c_{m-1}\right\}$ is a basis, for each $i=0, \ldots, p-1$ there is only one $S_{i} \in P(m)$ such that

$$
u_{i}=c_{S_{i}} .
$$

For each $S \in P(m)$, set

$$
r_{S}:=\prod_{\substack{\text { those } j \text { with } \\ S_{j}=S}} v_{j} .
$$

If no $S_{j}=S$, set $r_{S}=1$. Note that if there is an index $j$ with $S=S_{j}$, this index must be unique (because the expression $u_{i}=c_{S_{i}}$ is unique). Then by construction $r_{S} \in D_{G}\left(1, c_{S}\right)$ and in $k_{2}(G)$ we get

$$
\begin{aligned}
\sum_{k<m} \lambda\left(a_{k}\right) \lambda\left(b_{k}\right) & =\sum_{k<n} \lambda\left(c_{k}\right) \lambda\left(d_{k}\right)=\sum_{k<p} \lambda\left(u_{k}\right) \lambda\left(v_{k}\right) \\
& =\sum_{S \in P(m)} \lambda\left(\prod_{\substack{\text { those } \\
S_{j}=S}} c_{j}\right) \lambda\left(v_{j}\right) \\
& =\sum_{S \in P(m)} \sum_{\substack{\text { those j with } \\
S_{j}=S}} \lambda\left(c_{j}\right) \lambda\left(v_{j}\right) \\
& =\sum_{S \in P(m)} \lambda\left(c_{S}\right) \lambda\left(\prod_{\substack{\text { those } \\
S_{j}=S}} v_{S}\right) \\
& =\sum_{S \in P(m)} \lambda\left(c_{S}\right) \lambda\left(r_{S}\right)=\sum_{S \in P(m)} \sum_{k \in S} \lambda\left(c_{k}\right) \lambda\left(r_{S}\right) \\
& =\sum_{k<m} \lambda\left(c_{k}\right) \lambda\left(\prod_{S \in P(n)} r_{S}\right) .
\end{aligned}
$$

Since $\left\{c_{0}, \ldots, c_{m-1}\right\}$ is a basis, it follows that

$$
d_{i}=\prod_{S \in P(n)} r_{S}
$$

as desired.
(ii) $\Rightarrow$ (i). Under the hypothesis of (ii) we get

$$
\begin{aligned}
\sum_{k<n} \lambda\left(a_{k}\right) \lambda\left(b_{k}\right) & =\sum_{k<m} \lambda\left(c_{k}\right) \lambda\left(d_{k}\right)=\sum_{k<m} \lambda\left(c_{k}\right) \lambda\left(\prod_{S \in P(n)} r_{S}\right) \\
& =\sum_{k<m} \sum_{S \in P(m)} \lambda\left(c_{k}\right) \lambda\left(r_{S}\right)=\sum_{S \in P(m)} \sum_{k<m} \lambda\left(c_{i}\right) \lambda\left(r_{S}\right) \\
& =\sum_{S \in P(m)} \lambda\left(c_{S}\right) \lambda\left(r_{S}\right)=0 .
\end{aligned}
$$

Definition 4.2.4 (2.4 [30]).
$a-A$ reduced special group is [MC] if for all $n \leq 1$ and all form $\varphi$ over $G$,

$$
\text { For all } \sigma \in X_{G}, \text { if } \sigma(\varphi) \equiv 0 \bmod 2^{n} \text { then } \varphi \in I^{n} G \text {. }
$$

$b-A$ reduced special group is [SMC] if for all $n \geq 1$, the multiplication by $\lambda(-1)$ is an injection of $k_{n} G$ in $k_{n+1} G$.
An useful criteria for a reduced special group be [SMC] is given by:
Proposition 4.2.5 (2.5 [30]). Let $G$ be a reduced special group. Are equivalent:
$a-G$ is $S M C$;
$b$ - For all $n \geq 1, \varepsilon_{G}: k_{n} G \rightarrow B_{G}$ is injective.
Then, if $G$ is SMC, then $\varepsilon_{G}$ is an isomorphism between $k_{n} G$ and the subgroup $B_{G}(n)$ of $B_{G}$, for all $n \geq 1$.
2.6kt

Proposition 4.2.6 (2.6 [30]). Let $G$ be a formally real special group and $f: H \rightarrow G$ a complete embedding. If $H$ is [SMC], then $f_{*}$ is a graded ring monomorphism such that, for all $n \geq 0, f_{n}$ is injective.

An inductive system of special groups

$$
\mathcal{G}=\left(G_{i} ;\left\{f_{i j}: i \leq j \in I\right\}\right),
$$

provides an inductive system of graded ring, which nodes are $k_{*} G_{i}$ and morphisms are

$$
\left(f_{i j}\right)_{*}: k_{*} G_{i} \rightarrow k_{*} G_{j}, \text { for } i \leq j \text { in } I
$$

Theorem 4.2.7 (4.5 [30]). Let $\mathcal{G}=\left(G_{i} ;\left\{f_{i j}: i, j \in I, i \leq j\right\}\right)$ an inductive system of special groups over a directed poset $I$ and $\left(G ;\left\{f_{i}: i \in I\right\}\right)=\underline{\longrightarrow} \operatorname{G}$. Then $k_{*} G \cong \underset{i \in I}{\lim } k_{*} G_{i}$.

Corollary 4.2 .8 (4.6 [30]). The inductive limit of SMC groups is SMC.
If $S, T$ are $\mathbb{F}_{2}$-graded algebras with $S_{0}=T_{0}=\mathbb{F}_{2}$, the direct sum, $S \oplus T$, is the sequence of groups

$$
(S \oplus T)_{0}=\mathbb{F}_{2} \text { and }(S \oplus T)_{n}=S_{n} \oplus T_{n}, n \geq 1
$$

with the product defined by the rule $(x, y) \cdot(u, v)=(x u, y v)$. The $\mathbb{F}_{2}$-action on $S_{n} \oplus T_{n}$ is the usual action of $\mathbb{F}_{2}$-modules.
5.1kt

Theorem 4.2.9 (5.1 [30). Let $G_{1}, \ldots, G_{m}$ be special groups and $\prod_{i=1}^{m} G_{i}$. Then there exists a graded morphism

$$
\gamma: k_{*} P \rightarrow \bigoplus_{i=1}^{m} k_{*} G_{i},
$$

defined on the generators by the rule

$$
\gamma_{n}\left(\lambda\left(a_{1}\right) \ldots \lambda\left(a_{n}\right)\right)=\left\langle\lambda\left(\pi_{1}\left(a_{1}\right)\right) \ldots \lambda\left(\pi_{1}\left(a_{n}\right)\right) \ldots \lambda\left(\pi_{m}\left(a_{1}\right)\right) \ldots \lambda\left(\pi_{m}\left(a_{n}\right)\right)\right\rangle
$$

where $\pi_{i}: P \rightarrow G_{i}$ is the canonical projection, $i=1, \ldots, m$. Moreover, $\gamma$ send the multiplication by $\lambda(-1, \ldots,-1)$ on $P$ in the product $\lambda(-1) \ldots \lambda(-1)$ in $\bigoplus_{i=1}^{m} k_{*} G_{i}$.
5.4kt

Corollary 4.2.10 (5.4 [30]). The finite product of SMC groups is SMC.
5.6kt

Definition 4.2.11 (5.6 [30]). Let $\left\{G_{i}\right\}_{i \in I}$ be a family of special groups. Denote by $\bigoplus_{i \in I}^{*} G_{i}$ the following pre-special subgroup of $G=\prod_{i \in I} G_{i}$ :

$$
\bigoplus_{i \in I}^{*} G_{i}=\left\{x \in G: \text { exists } J \subseteq I \text { finite such that } x_{i}= \pm 1, \forall i \in I \backslash J\right\}
$$

with the special relation induced by the relation on $G$ and $-1=-1_{G}$. Such pre-special group will be called the $S G$-sum of the family $\left\{G_{i}\right\}_{i \in I}$.

In general, we do not have a canonical SG-embedding from $G$ into $G \times H$. On the other side, if we introduce a $\mathbb{Z}_{2}$ factor we can get around this situation. Let $I \subseteq J$ be finite sets and $G_{j}, j \in J$ be formally real special groups. Consider

$$
G_{J}:=\prod_{j \in J} G_{j}, G_{I}:=\prod_{i \in I} G_{i} .
$$

Let $\left\{G_{i}\right\}_{i \in I}$ be a family of formally real special groups. For each subset $A \subseteq I$, let $G_{A}=$ $\prod_{i \in A} G \times \mathbb{Z}_{2}$. If $A, B \subseteq I$ are finite subsets with $A \subseteq B$, we have a complete embedding $\alpha_{A B}$ : $G_{A} \rightarrow G_{B}$. Since the set $\mathcal{P}_{\text {Fin }}(I)$ of finite parts of $I$ with the inclusion order is up direct, we have the following inductive system of formally real special groups:

$$
\mathcal{G}=\left(G_{A} ;\left\{\alpha_{A B}: A, B \in \mathcal{P}_{F i n}(I), A \subseteq B\right\}\right)
$$

Theorem 4.2.12 (5.7 [30]). Let $\left\{G_{i}\right\}_{i \in I}$ be an infinite family of formally real special groups. Denote $S=\bigoplus_{i \in I}^{*} G_{i}$. With the above notations, we have $S=\underline{\longrightarrow} \mathcal{G}$. Moreover:
$a-k_{*} S=\underline{\longrightarrow} \mathcal{K} \mathcal{G}$, where $\mathcal{K}$ is the inductive system of the $K$-theory rings associated to $\mathcal{G}$;
b- The SG-sum of SMC groups is SMC.
6.8kt

Proposition 4.2.13 (6.8 [30]). Every extension of a SMC-group is SMC.
6.10kt

Theorem 4.2.14 (6.10 [30]). Let $G$ be a special group and $\Delta$ a group of exponent 2 finite or countable, of dimension $d \geq 1$ where considered as a $\mathbb{F}_{2}$-vector space. For each $n \geq 1$

$$
k_{n} G[\Delta]=\left\{\begin{array}{l}
\bigoplus_{j=0}^{d}\left(k_{n} G\right)^{\left.()^{d}\right)}, \text { if } d \text { is finite; } \\
k_{n} G \oplus\left[\bigoplus_{j=1}^{n}\left(\bigoplus_{d \geq 1} k_{n-j} G\right)\right], \text { if d is countable infinite. }
\end{array}\right.
$$

### 4.3 The K-theory for Multifields/Hyperfields

In this Section we introduce the notion of K-theory of a hyperfield essentially repeating the construction in 4.1.6 replacing the word "field" by "hyperfield" and explore some of this basic properties. In particular, Theorem 4.3 .8 is an extension of a result [59], that gives us some evidence, that apart from the obvious resemblance, more technical aspects of this new theory can be developed (but with other proofs) in multi-structure setting in parallel with classical K-theory.

Definition 4.3.1 (The K-theory of a Hyperfield). For a hyperfield $F, K_{*} F$ is the graded ring

$$
K_{*} F=\left(K_{0} F, K_{1} F, K_{2} F, \ldots\right)
$$

defined by the following rules: $K_{0} F:=\mathbb{Z} . K_{1} F$ is the multiplicative group $\dot{F}$ written additively. With this purpose, we fix the canonical "logarithm" isomorphism

$$
\rho: \dot{F} \rightarrow K_{1} F,
$$

where $\rho(a b)=\rho(a)+\rho(b)$. Then $K_{n} F$ is defined to be the quotient of the tensor algebra

$$
K_{1} F \otimes K_{1} F \otimes \ldots \otimes K_{1} F(n \text { times })
$$

by the (homogeneous) ideal generated by all $\rho(a) \otimes \rho(b)$, with $a \neq 0,1$ and $b \in 1-a$.
In other words, for each $n \geq 2$,

$$
K_{n} F:=T^{n}\left(K_{1} F\right) / Q^{n}\left(K_{1}(F)\right),
$$

where

$$
T^{n}\left(K_{1} F\right):=K_{1} F \otimes_{\mathbb{Z}} K_{1} F \otimes_{\mathbb{Z}} \ldots \otimes_{\mathbb{Z}} K_{1} F
$$

and $Q^{n}\left(K_{1}(F)\right)$ is the subgroup generated by all expressions of type $\rho\left(a_{1}\right) \otimes \rho\left(a_{2}\right) \otimes \ldots \otimes \rho\left(a_{n}\right)$ such that $a_{i} \in 1-a_{j}$ for some $i, j$ with $1 \leq i, j \leq n$.

To avoid carrying the over line symbol, we will adopt all the conventions used in DickmannMiraglia's K-theory (as explained in above Definition 4.2.1). Just as it happens with the previous K-theories, a generic element $\eta \in K_{n} F$ has the pattern

$$
\eta=\rho\left(a_{1}\right) \otimes \rho\left(a_{2}\right) \otimes \ldots \otimes \rho\left(a_{n}\right)
$$

for some $a_{1}, \ldots, a_{n} \in \dot{F}$, with $a_{i} \in 1-a_{i+1}$ for some $1 \leq i<n$. Note that if $F$ is a field, then " $b \in 1-a$ " just means $b=1-a$, and the hyperfield and Milnor's K-theory for $F$ coincide.

The very first task, is to extend the basic properties valid in Milnor's and Dickmann-Miraglia's K-theory to ours. Here we already need to restrict our attention to hyperbolic hyperfields:

Lemma 4.3.2 (Basic Properties I). Let $F$ be an hyperbolic hyperfield. Then
$a-\rho(1)=0$.
$b-$ For all $a \in \dot{F}, \rho(a) \rho(-a)=0$ in $K_{2} F$.
$c-$ For all $a, b \in \dot{F}, \rho(a) \rho(b)=-\rho(b) \rho(a)$ in $K_{2} F$.
$d$ - For every $a_{1}, \ldots, a_{n} \in \dot{F}$ and every permutation $\sigma \in S_{n}$,

$$
\rho\left(a_{1}\right) \ldots \rho\left(a_{i}\right) \ldots \rho\left(a_{n}\right)=\operatorname{sgn}(\sigma) \rho\left(a_{1}\right) \ldots \rho\left(a_{n}\right) \text { in } K_{n} F .
$$

$e-$ For every $\xi \in K_{m} F$ and $\eta \in K_{n} F, \eta \xi=(-1)^{m n} \xi \eta$ in $K_{m+n} F$.
$f$ - For all $a \in \dot{F}, \rho(a)^{2}=\rho(a) \rho(-1)$.
Proof.
a - Is an immediate consequence of the fact that $\rho$ is an isomorphism.
b- Since $F$ hyperbolic, $1-1=F$. Then $-a^{-1} \in 1-1$ for all $a \in \dot{F}$, and hence, $-1 \in-1+a^{-1}$. Multiplying this by $-a$, we get $a \in 1-a$. By definition, this imply $\rho(a) \rho(-a)=0$.
c-By item (b), $\rho(a b) \rho(-a b)=0$ in $K_{2} F$. But

$$
\begin{aligned}
\rho(a b) \rho(-a b) & =\rho(a) \rho((-a) b)+\rho(b) \rho((-b) a) \\
& =\rho(a) \rho(-a)+\rho(a) \rho(b)+\rho(b) \rho(-b)+\rho(b) \rho(a) \\
& =\rho(a) \rho(b)+\rho(b) \rho(a) .
\end{aligned}
$$

From $\rho(a) \rho(b)+\rho(b) \rho(a)=\rho(a b) \rho(-a b)=0$, we get the desired result $\rho(a) \rho(b)=-\rho(a) \rho(b)$ in $K_{2} F$.
d - This is a consequence of item (c) and an inductive argument.
e - This is a consequence of item (d) and an inductive argument, using the fact that an element in $K_{n} F$ has pattern

$$
\eta=\rho\left(a_{1}\right) \otimes \rho\left(a_{2}\right) \otimes \ldots \otimes \rho\left(a_{n}\right)
$$

for some $a_{1}, \ldots, a_{n} \in \dot{F}$, with $a_{i} \in 1-a_{j}$ for some $1 \leq i<j \leq n$.
f - Follow from the fact that $F$ is hyperbolic i.e, for all $a \in \dot{F}, a \in 1-1$.

An element $a \in \dot{F}$ induces a morphism of graded rings $\omega^{a}=\left\{\omega_{n}^{a}\right\}_{n \geq 1}: K_{*} F \rightarrow K_{*} F$ of degree 1 , where $\omega_{n}^{a}: K_{n} F \rightarrow K_{n+1} F$ is the multiplication by $\lambda(-a)$. When $a=-1$, we write

$$
\omega=\left\{\omega_{n}\right\}_{n \geq 1}=\left\{\omega_{n}^{-1}\right\}_{n \geq 1}=\omega^{-1} .
$$

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3.3ktmultiadap

Proposition 4.3.3 (Adapted from 3.3 of [30]). Let $F, K$ be hyperbolic hyperfields and $\varphi: F \rightarrow L$ be a morphism. Then $\varphi$ induces a morphism of graded rings

$$
\varphi_{*}=\left\{\varphi_{n}: n \geq 0\right\}: K_{*} F \rightarrow K_{*} L
$$

where $\varphi_{0}=I d_{\mathbb{Z}}$ and for all $n \geq 1, \varphi_{n}$ is given by the following rule on generators

$$
\varphi_{n}\left(\rho\left(a_{1}\right) \ldots \rho\left(a_{n}\right)\right)=\rho\left(\varphi\left(a_{1}\right)\right) \ldots \rho\left(\varphi\left(a_{n}\right)\right)
$$

Moreover if $\varphi$ is surjective then $\varphi_{*}$ is also surjective, and if $\psi: L \rightarrow M$ is another morphism then $a-(\psi \circ \varphi)_{*}=\psi_{*} \circ \varphi_{*}$ and $I d_{*}=I d$.
$b$ - For all $a \in \dot{F}$ the following diagram commute:

$c-$ If $\varphi(1)=1$ then for all $n \geq 1$ the following diagram commute:


Proof. Firstly, note that $\varphi$ extends to a function $\varphi_{1}: K_{1} F \rightarrow K_{1} L$ given by the rule

$$
\varphi_{1}(\rho(a))=\rho(\varphi(a))
$$

Certainly $\varphi_{1}$ is a morphism because

$$
\varphi_{1}(0)=\varphi_{1}(\rho(1))=\rho(\varphi(1))=\rho(1)=0
$$

and for all $\rho(a), \rho(b) \in K_{1} F$,

$$
\varphi_{1}(\rho(a)+\rho(b))=\varphi_{1}(\rho(a b))=\rho(\varphi(a b))=\rho(\varphi(a) \varphi(b))=\rho(\varphi(a))+\rho(\varphi(b))
$$

Proceeding inductively, for all $n \geq 1$ we extend $\varphi$ to a function $\varphi_{n}: \prod_{i=1}^{n} K_{1} F \rightarrow K_{n} L$ given by the rule

$$
\varphi\left(\rho\left(a_{1}\right), \ldots, \rho\left(a_{n}\right)\right):=\varphi_{1}\left(\rho\left(a_{1}\right)\right) \ldots \varphi_{1}\left(\rho\left(a_{n}\right)\right)=\rho\left(\varphi\left(a_{1}\right)\right) \ldots \rho\left(\varphi\left(a_{n}\right)\right)
$$

Then if $i=1, \ldots, n$ and $b_{i} \in k_{1} F$ we have

$$
\begin{aligned}
& \varphi_{n}\left(\rho\left(a_{1}\right), \ldots, \rho\left(a_{i}\right)+\rho\left(b_{i}\right), \ldots, \rho\left(a_{n}\right)\right)=\varphi_{n}\left(\rho\left(a_{1}\right), \ldots, \rho\left(a_{i} b_{i}\right), \ldots, \rho\left(a_{n}\right)\right)= \\
& \rho\left(\varphi\left(a_{1}\right)\right) \ldots \rho\left(\varphi\left(a_{i} b_{i}\right) \ldots \rho\left(\varphi\left(a_{n}\right)\right)=\rho\left(\varphi\left(a_{1}\right)\right) \ldots \rho\left(\varphi\left(a_{i}\right) \varphi\left(b_{i}\right)\right) \ldots \rho\left(\varphi\left(a_{n}\right)\right)=\right. \\
& \rho\left(\varphi\left(a_{1}\right) \ldots\left[\rho\left(\varphi\left(a_{i}\right)+\varphi\left(b_{i}\right)\right)\right] \ldots \rho\left(\varphi\left(a_{n}\right)\right)=\right. \\
& \rho\left(\varphi\left(a_{1}\right)\right) \ldots \rho\left(\varphi\left(a_{i}\right)\right) \ldots \rho\left(\varphi\left(a_{n}\right)\right)+\rho\left(\varphi\left(a_{1}\right)\right) \ldots \rho\left(\varphi\left(b_{i}\right)\right) \ldots \rho\left(\varphi\left(a_{n}\right)\right)= \\
& \varphi_{n}\left(\rho\left(a_{1}\right), \ldots, \rho\left(a_{i}\right), \ldots, \rho\left(a_{n}\right)\right)+\varphi_{n}\left(\rho\left(a_{1}\right), \ldots, \rho\left(b_{i}\right), \ldots, \rho\left(a_{n}\right)\right),
\end{aligned}
$$

then for each $n, \varphi_{n}: \prod_{i=1}^{n} K_{1} F \rightarrow K_{n} L$ is multilinear and by the universal property of tensor product there is an unique morphism

$$
\tilde{\varphi}_{n}: \bigotimes_{j=1}^{n} K_{1} F \rightarrow K_{n} L
$$

extending $\varphi_{n}$. By construction (and using the fact that $\varphi$ is a morphism), $\operatorname{Ker}\left(\tilde{\varphi}_{n}\right)=Q^{n}\left(K_{1} F\right)$, which provides an unique morphism $\bar{\varphi}_{n}: T^{n}\left(K_{1} F\right) / Q^{n}\left(K_{1}(F) \rightarrow K_{n} L\right.$ such that $\tilde{\varphi}_{n}=\bar{\varphi}_{n} \circ \pi_{n}$, where $\pi_{n}$ is the canonical projection $T^{n}\left(K_{1} F\right)$ in $Q^{n}\left(k_{1} F\right)$. Then taking $\varphi_{0}=I d_{\mathbb{Z}}$, we get a morphism $\varphi_{*}: K_{*} F \rightarrow K_{*} L$, given by $\varphi_{*}=\left\{\bar{\varphi}_{n}: n \geq 0\right\}$.

For items (a) and (b), it is enough to note that these properties holds for $\tilde{\varphi}_{n}, n \geq 0$, and after the application of projection, we get the validity for $\bar{\varphi}_{n}=\pi_{n} \circ \tilde{\varphi}_{n}$.

Item (c) follows by the same argument of items (a) and (b), noting that $\varphi(1)=1$ imply $\varphi(-1)=-1$. By abuse of notation, we denote

$$
\varphi_{*}=\left\{\bar{\varphi}_{n}: n \geq 0\right\}=\left\{\varphi_{n}: n \geq 0\right\} .
$$

We also have the reduced K-theory graded ring $k_{*} F=\left(k_{0} F, k_{1} F, \ldots, k_{n} F, \ldots\right)$ in the hyperfield context, which is defined by the rule $k_{n} F:=K_{n} F / 2 K_{n} F$ for all $n \geq 0$. Of course for all $n \geq 0$ we have an epimorphism $q: K_{n} F \rightarrow k_{n} F$ simply denoted by $q(a):=[a], a \in K_{n} F$. It is immediate that $k_{n} F$ is additively generated by $\left\{\left[\rho\left(a_{1}\right)\right] . .\left[\rho\left(a_{n}\right)\right]: a_{1}, \ldots, a_{n} \in \dot{F}\right\}$. We simply denote such a generator by $\tilde{\rho}\left(a_{1}\right) \ldots \tilde{\rho}\left(a_{n}\right)$ or even $\rho\left(a_{1}\right) \ldots \rho\left(a_{n}\right)$ whenever the context allows it.

We also have some basic properties of the reduced K-theory, which proof is just a translation of 2.1 of [30]:
2.1ktmulti

Lemma 4.3.4 (Adapted from 2.1 [30]). Let $F$ be a hyperbolic hyperfield, $x, y, a_{1}, \ldots, a_{n} \in \dot{F}$ and $\sigma$ be a permutation on $n$ elements.
$a$ - In $k_{2} F, \rho(a)^{2}=\rho(a) \rho(-1)$. Hence in $k_{m} F, \rho(a)^{m}=\rho(a) \rho(-1)^{m-1}, m \geq 2$;
$b-$ In $k_{2} F, \rho(a) \rho(-a)=\rho(a)^{2}=0$;
$c-\operatorname{In} k_{n} F, \rho\left(a_{1}\right) \rho\left(a_{2}\right) \ldots \rho\left(a_{n}\right)=\rho\left(a_{\sigma 1}\right) \rho\left(a_{\sigma 2}\right) \ldots \rho\left(a_{\sigma n}\right) ;$
$d-$ For $n \geq 1$ and $\xi \in k_{n} F, \xi^{2}=\rho(-1)^{n} \xi$;
$e$ - If $F$ is a real reduced hyperfield, then $x \in 1+y$ and $\rho(y) \rho\left(a_{1}\right) \ldots \rho\left(a_{n}\right)=0$ implies

$$
\rho(x) \rho\left(a_{1}\right) \rho\left(a_{2}\right) \ldots \rho\left(a_{n}\right)=0 .
$$

Moreover the results in Proposition 4.3.3 continue to hold if we took $\varphi_{*}=\left\{\varphi_{n}: n \geq 0\right\}$ : $k_{*} F \rightarrow k_{*} L$.
ktmarshall1
Proposition 4.3.5. Let $F$ be a hyperfield and $T \subseteq F$ be a multiplicative subset such that $F \subseteq T$. Then, for each $n \geq 1$

$$
K_{n}\left(F /{ }_{m} T^{*}\right) \cong k_{n}\left(F / m T^{*}\right) .
$$

Proof. Since $F^{2} \subseteq T$, for all $a \in F /{ }_{m} T^{*}$ we have

$$
0=\rho\left(a^{2}\right)=\rho(a)+\rho(a) .
$$

Then $2 K\left(F / m T^{*}\right)=0$ and we get $K_{n}\left(F /{ }_{m} T^{*}\right) \cong k_{n}\left(F / m T^{*}\right), n \geq 1$.
ktmarshall2
Theorem 4.3.6. Let $F$ be a hyperbolic hyperfield and $T \subseteq F$ be a multiplicative subset such that $F \subseteq T$. Then there is a surjective morphism

$$
k(F) \rightarrow k\left(F / m T^{*}\right) .
$$

Moreover, for each $n \geq 1$,

$$
k_{n}(F) \cong K_{n}\left(F / m \dot{F}^{2}\right) \cong k_{n}\left(F / m \dot{F}^{2}\right) .
$$

Before we prove it, we need a Lemma:
lemktmulti1
Lemma 4.3.7. Let $F$ be a hyperfield and $n \geq 1$. Then

$$
2 K_{n}(F)=\left\{\sum_{j=1}^{p} \rho\left(a_{j 1}\right) \ldots \rho\left(a_{j n}\right): \text { for all } j \text { there is an index } k \text { such that } a_{j k}=b_{i}^{2}, b_{i} \in \dot{F}\right\} .
$$

Proof. Let $\eta \in 2 K_{n} F$. Then

$$
\eta=\left(\sum_{j=1}^{p} \rho\left(a_{j 1}\right) \ldots \rho\left(a_{j n}\right)\right)+\left(\sum_{j=1}^{p} \rho\left(a_{j 1}\right) \ldots \rho\left(a_{j n}\right)\right), d_{i j} \in \dot{F} .
$$

By induction, we only need to consider the case $p=1$, so

$$
\rho\left(a_{1}\right) \ldots \rho\left(a_{n}\right)+\rho\left(a_{1}\right) \ldots \rho\left(a_{n}\right)=\rho\left(a_{1}^{2}\right) \rho\left(a_{2}\right) \ldots \rho\left(a_{n}\right) .
$$

and we get $\subseteq$. The reverse inclusion follow by the same calculation.

Proof of Theorem 4.3.6. Let $\pi: F \rightarrow F /{ }_{m} T^{*}$ denote the canonical projection. By Proposition 4.3 .3 there is a morphism $\pi_{*}: K(F) \rightarrow K\left(F /{ }_{m} T^{*}\right)$. Since $\pi$ is surjective, $\pi_{*}$ is surjective.

Now, let $\pi: F \rightarrow F /{ }_{m} \dot{F}^{2}$ and $q: K(F) \rightarrow k(F)$ the canonical projections. Denote elements in $F / m \dot{F}^{2}$ by $[a] \in F / m \dot{F}^{2}, a \in F$ and elements in $k_{n}(F)$ by $\tilde{\rho}\left(a_{1}\right) \ldots \tilde{\rho}\left(a_{n}\right)$. For all $n \geq 1$ we have an induced morphism $\tilde{q}_{n}: K_{n}\left(F /{ }_{m} \dot{F}^{2}\right) \rightarrow k_{n}(F)$ given by the rule

$$
\tilde{q}_{n}\left(\rho\left(\left[a_{1}\right]\right) \ldots \rho\left(\left[a_{n}\right]\right)\right):=\tilde{\rho}\left(a_{1}\right) \ldots \tilde{\rho}\left(a_{n}\right) .
$$

This morphism $\tilde{\pi}_{n}$ makes the following diagram commute

and then, $\tilde{q}_{n}$ is surjective. Finally, if $\tilde{q}_{n}\left(\rho\left(\left[a_{1}\right]\right) \ldots \rho\left(\left[a_{n}\right]\right)\right)=0$, then $\tilde{\rho}\left(a_{1}\right) \ldots \tilde{\rho}\left(a_{n}\right)=0$, and hence $\rho\left(a_{1}\right) \ldots \rho\left(a_{n}\right) \in 2 K_{n}(F)$. By Lemma 4.3.7

$$
\rho\left(a_{1}\right) \ldots \rho\left(a_{n}\right)=\sum_{j=1}^{p} \rho\left(d_{j 1}\right) \ldots \rho\left(d_{j n}\right), d_{i j} \in \dot{F}
$$

and for all $i$ there is an index $k$ such that $a_{i k}=b_{i}^{2}, b_{i} \in \dot{F}$. Therefore

$$
\begin{aligned}
\pi_{n}\left(\rho\left(a_{1}\right) \ldots \rho\left(a_{n}\right)\right) & =\pi_{n}\left(\sum_{j=1}^{p} \rho\left(d_{j 1}\right) \ldots \rho\left(d_{j n}\right)\right) \\
& \left.=\sum_{j=1}^{p} \pi_{n}\left(\rho\left(d_{j 1}\right) \ldots \rho\left(d_{j n}\right)\right)=\sum_{j=1}^{p} \rho\left(\left[d_{j 1}\right]\right) \ldots \rho\left(\left[d_{j n}\right)\right]\right) \\
& \left.=\sum_{j=1}^{p}\left[d_{j 1}\right]\right) \ldots \rho([1]) \ldots \rho\left(\left[d_{j n}\right)\right]=0 .
\end{aligned}
$$

Then $\operatorname{Ker}\left(\tilde{q}_{n}\right)=[0]$, proving that $\tilde{q}_{n}$ is injective. Then $\tilde{q}_{n}$ is an isomorphism, and composing all the isomorphisms obtained here we get

$$
k(F) \cong K\left(F /{ }_{m} \dot{F}^{2}\right) \cong k\left(F /{ }_{m} \dot{F}^{2}\right) .
$$

The Theorem 4.3 .8 below generalizes Proposition 5.10 of [59]: this constitutes a fundamental technical step to build profinite (Galois) groups associated to a pre-special hyperfield in [20].

Lets establish some notation: for $n \geq 0$ we denote

$$
P(n)=\mathcal{P}(\{0, \ldots, n-1\}) \backslash\{\emptyset\}
$$

and for $0 \leq i \leq n-1$, denote

$$
P(n, i)=\{X \in P(n): i \in X\} .
$$

For a be a pre-special hyperfield $F$ and $\left\{a_{0}, \ldots, a_{n-1}\right\} \subseteq F^{*} \mathbb{F}_{2}$-linearly independent, if $S \in P(n)$ we denote

$$
a_{S}:=a_{0}^{\varepsilon_{0}} \ldots a_{n-1}^{\varepsilon_{n-1}},
$$

where $\varepsilon_{0} \in\{0,1\}$ for all $i=0, . ., n-1$ and $\varepsilon_{i}=1$ if and only if $i \in S$.

Remember that by the very Definition of $k_{n}(F)$,

$$
k_{2}(F):=\left[k_{1}(G) \otimes k_{1}(G)\right] / M,
$$

where $M$ is the subgroup of $k_{1}(G) \otimes k_{1}(G)$ generated by

$$
\left\{\rho(a) \rho(b): a \in D_{G}(1,-b)\right\} .
$$

Theorem 4.3.8. Let $F$ be a pre-special hyperfield and $\left\{a_{0}, \ldots, a_{n-1}\right\} \subseteq F^{*} \mathbb{F}_{2}$-linearly independent. The following conditions are equivalent:
$i$ - There exists $\left\{b_{0}, \ldots, b_{n-1}\right\} \subseteq F^{*}$ such that

$$
\sum_{k<n} \rho\left(a_{k}\right) \rho\left(b_{k}\right)=0 \text { in } k_{2}(F) .
$$

ii - There exist subsets $\left\{c_{0}, \ldots, c_{m-1}\right\},\left\{d_{0}, \ldots, d_{n-1}\right\}$ of $F^{*}$ with $m \geq n$ such that
(a) $\left\{c_{0}, \ldots, c_{m-1}\right\}$ is linearly independent and $c_{i}=a_{i}$ for all $i<n$;
(b) $d_{i}=b_{i}$ for all $i<n$ and $d_{i}=1$ for $i=n, \ldots, m-1$.
(c) For all $x \in C:=\left[c_{0}, \ldots, c_{m-1}\right]$, there is some $r_{x} \in(1-x) \backslash\{0\}$ such that for each $i<m$

$$
d_{i}=\prod_{x \in C_{i}} r_{x}
$$

where

$$
C_{i}=\left\{\prod_{k<m} c_{k}^{\varepsilon_{k}}: \varepsilon_{k} \in\{0,1\} \text { and } \varepsilon_{i}=1\right\} .
$$

In other words, $C_{i}$ is "counting" all products $c_{0}^{r_{0}} \ldots c_{i}^{1} \ldots c_{m-1}^{r_{m-1}}$. Since for all $x \in C:=\left[c_{0}, \ldots, c_{m-1}\right]$ there exist $S \in P(m)$ such that

$$
x=\prod_{i \in S} c_{i}:=c_{S}
$$

Denoting $r_{x}$ by $r_{S}$ we can rewrite

$$
d_{i}=\prod_{x \in C_{i}} r_{x}=\prod_{S \in P(m)} r_{S} .
$$

Proof of Theorem 4.3.8. (i) $\Rightarrow$ (ii). Let

$$
\sum_{k<n} \rho\left(a_{k}\right) \rho\left(b_{k}\right)=0 \text { in } k_{2}(F) .
$$

Then there exist $u_{0}, \ldots, u_{p-1}, v_{0}, \ldots, v_{p-1} \in F^{*}$ such that $v_{i} \in 1+u_{i}$ for $i=0, \ldots, p-1$ and

$$
\sum_{k<n} \rho\left(a_{k}\right) \rho\left(b_{k}\right)=\sum_{k<n} \rho\left(a_{k}\right) \rho\left(b_{k}\right) \text { in } k_{1}(F) \otimes k_{1}(F) .
$$

Enlarge the set $\left\{a_{0}, \ldots, a_{n-1}\right\}$ to a base $\left\{c_{0}, \ldots, c_{m-1}\right\}$ of $\left[\left\{a_{0}, \ldots, a_{n-1}, u_{0}, \ldots, u_{p-1}\right\}\right]$, with $c_{i}=a_{i}$ for
all $i<n$. For all $x \in C:=\left[c_{0}, \ldots, c_{m-1}\right]$ there exist $S \in P(m)$ such that

$$
x=\prod_{i \in S} c_{i}:=c_{S} .
$$

Moreover, since $\left\{c_{0}, \ldots, c_{m-1}\right\}$ is a basis, for each $i=0, \ldots, p-1$ there is only one $S_{i} \in P(m)$ such that $u_{i}=c_{S_{i}}$. For each $S \in P(m)$, set

$$
r_{S}:=\prod_{\substack{\text { those } \text { with } \\ S_{j}=S}} v_{j} .
$$

If no $S_{j}=S$, set $r_{S}=1$. Note that if there is an index $j$ with $S=S_{j}$, this index must be unique (because the expression $u_{i}=c_{S_{i}}$ is unique). Then by construction $r_{S} \in 1+c_{S} \backslash\{0\}$ and in $k_{2}(F)$ we get

$$
\begin{aligned}
\sum_{k<m} \rho\left(a_{k}\right) \rho\left(b_{k}\right) & =\sum_{k<n} \rho\left(c_{k}\right) \rho\left(d_{k}\right)=\sum_{k<p} \rho\left(u_{k}\right) \rho\left(v_{k}\right) \\
& =\sum_{S \in P(m)} \rho\left(\prod_{\substack{\text { those jwith } \\
S_{j}=S}} c_{j}\right) \rho\left(v_{j}\right) \\
& =\sum_{S \in P(m)} \sum_{\substack{\text { those j with } \\
S_{j}=S}} \rho\left(c_{j}\right) \rho\left(v_{j}\right) \\
& =\sum_{S \in P(m)} \rho\left(c_{S}\right) \rho\left(\prod_{\substack{\text { those } j w i t h \\
S_{j}=S}} v_{S}\right) \\
& =\sum_{S \in P(m)} \rho\left(c_{S}\right) \rho\left(r_{S}\right)=\sum_{S \in P(m)} \sum_{k \in S} \rho\left(c_{k}\right) \rho\left(r_{S}\right) \\
& =\sum_{k<m} \rho\left(c_{k}\right) \rho\left(\prod_{S \in P(n)} r_{S}\right) .
\end{aligned}
$$

Since $\left\{c_{0}, \ldots, c_{m-1}\right\}$ is a basis, it follows that $d_{i}=\prod_{S \in P(n)} r_{S}$ as desired.
(ii) $\Rightarrow$ (i). Under the hypotheses of (ii) we get

$$
\begin{aligned}
\sum_{k<n} \rho\left(a_{k}\right) \rho\left(b_{k}\right) & =\sum_{k<m} \rho\left(c_{k}\right) \rho\left(d_{k}\right)=\sum_{k<m} \rho\left(c_{k}\right) \rho\left(\prod_{S \in P(n)} r_{S}\right) \\
& =\sum_{k<m} \sum_{S \in P(m)} \rho\left(c_{k}\right) \rho\left(r_{S}\right)=\sum_{S \in P(m)} \sum_{k<m} \rho\left(c_{i}\right) \rho\left(r_{S}\right) \\
& =\sum_{S \in P(m)} \rho\left(c_{S}\right) \rho\left(r_{S}\right)=0 .
\end{aligned}
$$

### 4.4 Inductive Graded Rings: An Abstract Approach

After the three K-theories defined in the above sections, it is desirable (or, at least, suggestive) the rise of an abstract environment that encapsule all them, and of course, provide an axiomatic approach to guide new extensions of the concept of K-theory in the context of the algebraic and abstract theories of quadratic forms. The inductive graded rings fits this purpose. Here we will present three versions. The first one is:
igr1
Definition 4.4.1 (Inductive Graded Rings First Version (adapted from Definition 9.7 of [28])). An inductive graded ring (or Igr for short) is a structure $R=\left(\left(R_{n}\right)_{n \geq 0},\left(h_{n}\right)_{n \geq 0}, *_{n m}\right)$ where

$$
i-R_{0} \cong \mathbb{F}_{2} .
$$

ii - $R_{n}$ has a group structure $\left(R_{n},+, 0, \top_{n}\right)$ of exponent 2 with a distinguished element $T_{n}$.
iii - $h_{n}: R_{n} \rightarrow R_{n+1}$ is a group homomorphism such that $h_{n}\left(\top_{n}\right)=\top_{n+1}$.
iv - For all $n \geq 1, h_{n}=*_{1 n}\left(\top_{1},{ }_{-}\right)$.
$v$ - The binary operations $*_{n m}: R_{n} \times R_{m} \rightarrow R_{n+m}, n, m \in \mathbb{N}$ induces a commutative ring structure on the abelian group

$$
R=\bigoplus_{n \geq 0} R_{n}
$$

with $1=\mathrm{T}_{0}$.
vi- For $0 \leq s \leq t$ define

$$
h_{s}^{t}=\left\{\begin{array}{l}
I d_{R_{s}} \text { if } s=t \\
h_{t-1} \circ \ldots \circ h_{s+1} \circ h_{s} \text { if } s<t .
\end{array}\right.
$$

Then if $p \geq n$ and $q \geq m$, for all $x \in R_{n}$ and $y \in R_{m}$,

$$
h_{n}^{p}(x) * h_{m}^{q}(y)=h_{n+m}^{p+q}(x * y) .
$$

A morphism between Igr's $R$ and $S$ is a pair $f=\left(f,\left(f_{n}\right)_{n \geq 0}\right)$ where $f_{n}: R_{n} \rightarrow S_{n}$ is a morphism of pointed groups and

$$
f=\bigoplus_{n \geq 0} f_{n}: R \rightarrow S
$$

is a morphism of commutative rings with unity. The category of inductive graded rings (in first version) and their morphisms will be denoted by Igr.

A first consequence of these definitions is that: if

$$
f:\left(\left(R_{n}\right)_{n \geq 0},\left(h_{n}\right)_{n \geq 0}, *_{n m}\right) \rightarrow\left(\left(S_{n}\right)_{n \geq 0},\left(l_{n}\right)_{n \geq 0}, *_{n m}\right)
$$

is a morphism of Igr's then $f_{n+1} \circ h_{n}=l_{n} \circ f_{n}$.


In fact, since $R_{0} \cong \mathbb{F}_{2} \cong S_{0}$ and $f(1)=1$, then $f_{0}: R_{0} \rightarrow S_{0}$ is the unique abelian group isomorphism and $f_{1} \circ h_{0}=l_{0} \circ f_{0}$. If $n \geq 1$, for all $a_{n} \in R_{n}$ holds

$$
\begin{aligned}
f_{n+1} \circ h_{n}\left(a_{n}\right) & =f_{n+1} \circ\left(*_{1 n}\left(\top_{1}, a_{n}\right)\right)=f_{1}\left(\top_{1}\right) *_{1 n} f_{n}\left(a_{n}\right) \\
& =\top_{1} *_{1 n} f_{n}\left(a_{n}\right)=l_{n}\left(f_{n}\left(a_{n}\right)\right)=l_{n} \circ f_{n}\left(a_{n}\right)
\end{aligned}
$$

## Example 4.4.2.

$a$ - Let $F$ be a field of characteristic not 2. The main actors here are WF, the Witt ring of $F$ and $I F$, the fundamental ideal of $W F$. Is well know that $I^{n} F$, the $n$-th power of $I F$ is additively generated by $n$-fold Pfister forms over $F$. Now, let $R_{0}=W F / I F \cong \mathbb{F}_{2}$ and $R_{n}=I^{n} F / I^{n+1} F$. Finally, let $h_{n}=_{-} \otimes\langle 1,1\rangle, n \geq 1$. With these prescriptions we have an inductive graded ring $R$ associated to $F: W_{*}(F)$, the graded Witt ring of the field $F$.
$b$ - The previous example still works if we change the Witt ring of a field $F$ for the Witt ring of a (formally real) special group $G$.
$c$ - An inductive graded ring can be seen as a graded $F_{2}$-algebra $R$ with $R_{0}=F_{2}$ and a distinguished element $T_{1}$ in $R_{1}$.

There is an alternative definition for Igr with a first-order theoretic flavor. It is a technical framework that allows achieving some model-theoretic results.

Before define it, we need some preparation. First of all, we set up the language. Here, we will work with the poli-sorted framework (as established in chapter 5 of [1]), which means the following:

Let $S$ be a set (of sorts). For each $s \in S$ assume a countable set $\operatorname{Var}_{s}$ of variables of sort $s$ (with the convention if $s \neq t$ then $\operatorname{Var}_{s} \cap \operatorname{Var}_{t}=\emptyset$ ). For each sort $s \in S$ an equality symbol $={ }_{s}$ (or just $=$ ); the connectives $\neg, \wedge, \vee, \rightarrow$ (not, and, or, implies); the quantifiers $\forall, \exists$ (for all, there exists).

A finitary $S$-sorted language (or signature) is a set $\mathcal{L}=(\mathcal{C}, \mathcal{F}, \mathcal{R})$ where:
i - $\mathcal{C}$ is the set of constant symbols. For each $c \in \mathcal{C}$ we assign an element $s \in S$, the sort of $c$;
ii - $\mathcal{F}$ is the set of functional symbols. For each $f \in \mathcal{F}$ we assign elements $s, s_{1}, \ldots, s_{n} \in S$, we say that $f$ has arity $s_{1} \times \ldots \times s_{n}$ and $s$ is the value sort of $f$; and we use the notation $f: s_{1} \times \ldots \times s_{n} \rightarrow s$.
iii - $\mathcal{R}$ is the set of relation symbols. $c \in \mathcal{C}$ we assign elements $s_{1}, \ldots, s_{n} \in S$, the arity of $R$; and we say that $R$ has arity $s_{1} \times \ldots \times s_{n}$.

A $\mathcal{L}$-structure $\mathcal{M}$ is, in this sense, prescribed by the following data:
i- The domain or universe of $\mathcal{M}$, which is an $S$-sorted set $|\mathcal{M}|:=\left(M_{s}\right)_{s \in S}$.
ii- For each constant symbol $c \in \mathcal{C}$ of arity $s$, an element $c^{\mathcal{M}} \in M_{s}$.
iii- For each functional symbol $f \in \mathcal{F}, f: s_{1} \times \ldots \times s_{n} \rightarrow s$, a function $f^{\mathcal{M}}: M_{s_{1}} \times \ldots \times M_{s_{n}} \rightarrow M_{s}$.
iv- For each relation symbol $R \in \mathcal{R}$ of arity $s_{1} \times \ldots \times s_{n}$ a relation, i.e. a subset $R^{\mathcal{M}} \subseteq M_{s_{1}} \times \ldots \times M_{s_{n}}$.
A $\mathcal{L}$-morphism $\varphi: \mathcal{M} \rightarrow \mathcal{N}$ is a sequence of functions $\varphi=\left(\varphi_{s}\right)_{s}:|\mathcal{M}| \rightarrow|\mathcal{N}|$ such that
i - for all $c \in \mathcal{C}$ of arity $s, \varphi_{s}\left(c^{\mathcal{M}}\right)=c^{\mathcal{N}} ;$
ii - for all $f: s_{1} \times \ldots \times s_{n} \rightarrow s$, if $\left(a_{1}, \ldots, a_{n}\right) \in: M_{s_{1}} \times \ldots \times M_{s_{n}}$, then $\varphi_{s}\left(f^{\mathcal{M}}\left(a_{1}, \ldots, a_{n}\right)\right)=$ $f^{\mathcal{N}}\left(\varphi_{s_{1}}\left(a_{1}\right), \ldots, \varphi_{s_{n}}\left(a_{n}\right)\right) ;$
iii - for all $R$ of arity $s_{1} \times \ldots \times s_{n}$, if $\left(a_{1}, \ldots, a_{n}\right) \in R^{\mathcal{M}}$ then $\left(\varphi\left(a_{1}\right), \ldots, \varphi\left(a_{n}\right)\right) \in R^{\mathcal{N}}$.
The category of $\mathcal{L}$-structures and $\mathcal{L}$-morphism in the poli-sorted language $\mathcal{L}$ will be denoted by $\operatorname{Str}_{s}(\mathcal{L})$.

The terms, formulas, occurrence and free variables definitions for the poli-sorted case are similar to the usual (single-sorted) first order ones. For example, the terms are defined as follows:
i - variables $x \in \operatorname{Var}_{s}$ and constants $c \in C_{s}$ are terms of value sort $s$;
ii - if $\vec{s}=\left\langle s_{1}, \ldots, s_{n}, s\right\rangle \in S^{n+1}, f \in \mathcal{F}$ with $f: s_{1} \times \ldots \times s_{n} \rightarrow s$, and $\tau_{1}, \ldots, \tau_{n}$ are terms of value sorts $s_{1}, \ldots, s_{n}$ respectively, then $f\left(\tau_{1}, \ldots, \tau_{n}\right)$ is a term of sort $s$.

As usual, we may write $\tau: s$ to indicate that the term $\tau$ has value sort $s$.
For the formulas:
i - if $x, y \in \operatorname{Var}_{s}$ then $x=y$ is a formula; if $\vec{s}=\left\langle s_{1}, \ldots, s_{n}\right\rangle \in S^{n}, R \in \mathcal{R}$ of arity $s_{1} \times \ldots \times s_{n}$ and $\tau_{1}, \ldots, \tau_{n}$ are terms of sort $s_{1}, \ldots, s_{n}$ respectively, then $R\left(\tau_{1}, \ldots, \tau_{n}\right)$ is a formula. These are the atomic formulas.
ii - If $\varphi_{1}, \varphi_{2}$ are formulas, then $\neg \varphi_{1}, \varphi_{1} \wedge \varphi_{2}, \varphi_{1} \vee \varphi_{2}$ and $\varphi_{1} \rightarrow \varphi_{2}$ are formulas.
iii - If $\varphi$ is a formula and $x \in \operatorname{Var}_{s}(s \in S)$, then $\forall x \varphi$ and $\exists x \varphi$ are formulas.
In our particular case, the set of sorts will be just $\mathbb{N}$. Then, for each $n, m \geq 0$, we set the following data:
i - $0_{n}, T_{n}$ are constant symbols of arity $n$. We use $0_{0}=0$ and $T_{0}=1$.
ii - $+_{n}: n \times n \rightarrow n$ is a binary operation symbol.
iii - $h_{n}: n \rightarrow(n+1)$ and $*_{n, m}: n \times m \rightarrow(n+m)$ are functional symbols.
The (first order) language of inductive graded rings $\mathcal{L}_{i g r}$ is just the following language (in the poli-sorted sense):

$$
\mathcal{L}_{i g r}:=\left\{0_{n}, \top_{n},+_{n}, h_{n}, *_{n m}: n, m \geq 0\right\} .
$$

The (first order) theory of inductive graded rings $T\left(\mathcal{L}_{i g r}\right)$ is the $\mathcal{L}_{i g r}$-theory axiomatized by the following $\mathcal{L}_{i g r}$-sentences, where we use $\cdot_{n}: 0 \times n \rightarrow n$ as an abbreviation for ${ }^{0}{ }_{n}$ :
i - For $n \geq 0$, sentences saying that " $+_{n}, 0_{n}, \top_{n}$ induces a pointed left $\mathbb{F}_{2}$-module":

$$
\begin{aligned}
& \forall x: n \forall y: n \forall z: n\left(\left(x+{ }_{n} y\right)+_{n} z=x+{ }_{n}\left(y+{ }_{n} z\right)\right) \\
& \forall x: n\left(x+{ }_{n} 0_{n}=x\right) \\
& \forall x: n \forall y: n\left(x+_{n} y=y+{ }_{n} x\right) \\
& \forall x: n\left(x+{ }_{n} x=0_{n}\right) \\
& \forall x: n(1 \cdot n x=x) \\
& \forall x: n \forall y: n \forall a: 0\left(a \cdot{ }_{n}\left(x+_{n} y\right)=a \cdot n x+_{n} a \cdot{ }_{n} y\right) \\
& \forall x: n \forall a: 0 \forall b: 0\left(\left(a+{ }_{0} b\right) \cdot{ }_{n} x=a \cdot_{n} x+_{n} b \cdot{ }_{n} x\right)
\end{aligned}
$$

ii - For $n \geq 0$, sentences saying that " $h_{n}$ is a pointed $\mathbb{F}_{2}$-morphism":

$$
\begin{aligned}
& \forall x: n \forall y: n\left(h_{n}\left(x+_{n} y\right)=h_{n}(x)+_{n+1} h_{n}(y)\right) \\
& \forall x: n \forall a: 0\left(h_{n}\left(a \cdot{ }_{n} x\right)=a \cdot{ }_{n} h_{n}(x)\right) \\
& h_{n}\left(\top_{n}\right)=\top_{n+1}
\end{aligned}
$$

iii - Sentences saying that " $R_{0} \cong \mathbb{F}_{2}$ ":

$$
\begin{aligned}
& 0_{0} \neq \top_{0} \\
& \forall x: n\left(x=0_{0} \vee x=\top_{0}\right)
\end{aligned}
$$

iv - Using the abbreviation $*_{n, m}(x, y)=x *_{n, m} y$, we write for $n, m \geq 0$ sentences saying that " $*_{n, m}$ is a biadditive function compatible with $h_{n}$ ":

$$
\begin{aligned}
& \forall x: n \forall y: n \forall z: m\left(\left(\left(x+_{n} y\right) *_{n m} z\right)=\left(x *_{m n} z+_{n+m} y *_{n m} z\right)\right) \\
& \forall x: n \forall y: m \forall z: m\left(\left(x *_{m n}\left(y+_{m} z\right)\right)=\left(x *_{n m} y+_{n+m} x *_{n m} z\right)\right) \\
& \forall x: n \forall y: m\left(h_{n+m}\left(x *_{n m} y\right)=h_{n}(x) *_{n m} h_{m}(y)\right)
\end{aligned}
$$

v - Sentences describing "the induced ring with product induced by $*_{n, m}, n, m \geq 0$ ":

$$
\begin{aligned}
& \forall x: n \forall y: m \forall z: p\left(\left(x *_{n, m} y\right) *_{(m+n), p} z=x *_{n,(m+p)}\left(y *_{m, p} z\right)\right) \\
& \forall x: n \forall y: m\left(x *_{n, m} y=y *_{m, n} x\right)
\end{aligned}
$$

vi - For $n \geq 1$, sentences saying that " $h_{n}=T_{1} *_{1 n-}$ ":

$$
\forall x: n\left(h_{n}(x)=\top_{1} *_{1 n} x\right)
$$

Now we are in a position to define another version of Igr:
Definition 4.4.3 (Inductive Graded Rings Second Version). An inductive graded ring (or (Igr) for short) is a model for $T\left(\mathcal{L}_{i g r}\right)$, or in other words, a $\mathcal{L}_{\text {igr }}$-structure $\mathcal{R}$ such that $\mathcal{R} \models \mathcal{L}_{\text {igr }} T\left(\mathcal{L}_{\text {igr }}\right)$. We denote by $I g r_{2}$ the category of $\mathcal{L}_{i g r}$-structures and $\mathcal{L}_{i g r}$-morphisms.

Again, after some straightforward calculations we can check:
Theorem 4.4.4. The categories Igr, Igr $_{2}$ are equivalent.
igr-re
Remark 4.4.5. Following a well-known procedure, it is possible to correspond theories on polysorted first-order languages with theories on traditional (single-sorted) first-order languages in such a way that the corresponding categories of models are equivalent. This allows a useful interchanging between model-theoretic results, in both directions. In particular, in the following, we will freely interchange the three notions of Igr indicated in this section.

### 4.5 Interchanging K-theories

We finalize this chapter with an use of Igr's to interchanging the three K-theory notions presented before in a functorial fashion. Lets first, look more carefully at theorem 4.5.1. We make the
following distinctions between K-theories:
$K$ will denote the Milnor's K-Theory,
$K^{d m}$ will denote the Dickmann-Miraglia's K-Theory,
$K^{\text {mult }}$ will denote the K-Theory of Hyperfields.

Of course, we need the following theorem:

## Theorem 4.5.1.

$a$ - Let $F$ be a field. Then $k_{*}^{m i l} F$ (the reduced Milnor $K$-theory) is an inductive graded ring.
$b$ - Let $G$ be a special group. Then $k_{*}^{d m} G$ (the Dickmann-Miraglia K-theory of $G$ ) is an inductive graded ring.
$c$ - Let $F$ be a hyperfield. Then $k_{*}^{\text {mult }} F$ (our reduced $K$-theory) is an inductive graded ring.
Proof. Item (a) is the content of Lemma 9.11 in [28], and item (b) is the content of Lemma 9.12 in [28. We prove item (c) and items (a) and (b) will proceed by the same argument.

Let $k_{*}^{\text {mult }} F=\left(k_{0} F, k_{1} F, \ldots, k_{n} F, \ldots\right)$ be the reduced K-theory of a hyperfield $F$. Let $\mathrm{T}_{0}=1$ and for each $n \geq 1$, we set $\top_{n}=l(-1)^{n}$ as the distinguished element of $m-n$. For each $n \geq 0$, let $\theta_{n}: \prod_{j=1}^{n} K_{1}^{\text {mult }} F \rightarrow \otimes_{j=1}^{n+1} K_{n+1}^{\text {mult }} F$ given by the rule

$$
\theta_{n}\left(\rho\left(a_{1}\right), \ldots, \rho\left(a_{n}\right)\right):=\rho(-1) \rho\left(a_{1}\right) \ldots \rho\left(a_{n}\right) .
$$

We have for each $i \in\{1, \ldots, n\}$ and each $a_{1}, \ldots, a_{n}, b_{i} \in F^{*}$ that

$$
\begin{aligned}
& \theta_{n}\left(\rho\left(a_{1}\right), \ldots \rho\left(a_{i}\right)+\rho\left(b_{i}\right), \ldots, \rho\left(a_{n}\right)\right)=\theta_{n}\left(\rho\left(a_{1}\right), \ldots \rho\left(a_{i} b_{i}\right), \ldots, \rho\left(a_{n}\right)\right):= \\
& \rho(-1) \rho\left(a_{1}\right) \ldots \rho\left(a_{i} b_{i}\right) \ldots \rho\left(a_{n}\right)=\rho(-1) \rho\left(a_{1}\right) \ldots\left[\rho\left(a_{i}\right)+\rho\left(b_{i}\right)\right] \ldots \rho\left(a_{n}\right)= \\
& \rho(-1) \rho\left(a_{1}\right) \ldots \rho\left(a_{i}\right) \ldots \rho\left(a_{n}\right)+\rho(-1) \rho\left(a_{1}\right) \ldots \rho\left(b_{i}\right) \ldots \rho\left(a_{n}\right)= \\
& \theta_{n}\left(\rho\left(a_{1}\right), \ldots, \rho\left(a_{i}\right), \ldots, \rho\left(a_{n}\right)\right)+\theta_{n}\left(\rho\left(a_{1}\right) \ldots \rho\left(b_{i}\right) \ldots \rho\left(a_{n}\right)\right),
\end{aligned}
$$

then $\theta_{n}$ is multilinear. By the universal property of tensor product, we have a group homomorphism $\tilde{\theta}_{n}: K_{n}^{\text {mult }} F \rightarrow K_{n+1}^{\text {mult }} F$ given by the rule ${ }^{2}$

$$
\tilde{\theta}_{n}\left(\rho\left(a_{1}\right) \ldots \rho\left(a_{n}\right)\right)=\rho(-1) \rho\left(a_{1}\right) \ldots \rho\left(a_{n}\right) .
$$

In order to make distinctions between reduced and non-reduced K-theories, we punctually denote an element in $k_{n}^{\text {mult }} F:=K_{n}^{\text {mult }} F / 2 K_{n}^{\text {mult }} F$ by $\tilde{\rho}\left(a_{1}\right) \ldots \tilde{\rho}\left(a_{n}\right)$. Lets also denote the canonical projection by $\pi_{n}: K_{n}^{\text {mult }} F \rightarrow k_{n}^{\text {mult }} F$. We define $\omega_{n}: k_{n}^{\text {mult }} F \rightarrow k_{n+1}^{\text {mult }} F$ by the following rule (on generators): for $a_{1}, \ldots, a_{n} \in \dot{F}$,

$$
\omega_{n}\left(\tilde{\rho}\left(a_{1}\right) \ldots \tilde{\rho}\left(a_{n}\right)\right)=\tilde{\rho}(-1) \tilde{\rho}\left(a_{1}\right) \ldots \tilde{\rho}\left(a_{n}\right) .
$$

In fact, if $\rho\left(a_{1}\right) \ldots \rho\left(b_{n}\right)-\rho\left(b_{1}\right) \ldots \rho\left(b_{n}\right) \in 2 K_{n}^{\text {mult }} F$ then

$$
\rho(-1) \rho\left(a_{1}\right) \ldots \rho\left(b_{n}\right)-\rho(-1) \rho\left(b_{1}\right) \ldots \rho\left(b_{n}\right)=\rho(-1)\left[\rho\left(a_{1}\right) \ldots \rho\left(b_{n}\right)-\rho\left(b_{1}\right) \ldots \rho\left(b_{n}\right)\right] \in 2 K_{n+1}^{\text {mult }} F,
$$

[^13]which proves that $\omega_{n}$ is in fact a group homomorphism making the following diagram commute


With these rules, we already have the properties (i)-(iv) of Definition 4.4.1 holding in $k_{*} F$, remaining only property (v). Note that $\omega_{s}^{t}=\tilde{\rho}(-1)^{t-s}$ for $0 \leq s<t$.

Now let $m, n, p, q \in \mathbb{N}, p \geq n, q \geq m$ and consider $x \in k_{n}^{m u l t} F$ and $y \in k_{m}^{m u l t} F$. Note that $\omega_{n}^{p}(x)=\tilde{\rho}(-1)^{p-n} \cdot x$ and $\omega_{m}^{q}(y)=\tilde{\rho}(-1)^{q-m} \cdot y$. Then

$$
\omega_{n}^{p}(x) \cdot \omega_{m}^{q}(y)=\left(\tilde{\rho}(-1)^{p-n} \cdot x\right)\left(\tilde{\rho}(-1)^{q-m} \cdot y\right)=\tilde{\rho}(-1)^{p-n+q-m} \cdot(x \cdot y)=\omega_{n-m}^{p+q}(x),
$$

completing the proof.
Using this Theorem (in addition with the argument of Lemma 3.3 in [30]) we obtain the following.
km2
Corollary 4.5.2. We have a functor and $k:$ Field $_{2} \rightarrow$ Igr induced by K-theory and Milnor's reduced $K$-theory.

Corollary 4.5.3. We have a functor $k^{\text {mult }}:$ MField $_{2} \rightarrow$ Igr induced by our reduced $K$-theory.
km4
Theorem 4.5.4 (Theorem 2.5 in [29]). Let $F$ be a field. The functor $G:$ Field $_{2} \rightarrow S G$ provides a functor $k_{*}^{d m}:$ Field $_{2} \rightarrow$ Igr (the special group K-theory functor) given on the objects by $k_{*}^{d m}(F)$ : $k_{*}^{d m}(G(F))$ and on the morphisms $f: F \rightarrow K$ by $k_{*}^{d m}(f):=G(f)_{*}$ (in the sense of Lemma 3.3 of [30]). Moreover, this functor commutes with the functors $G$ and $k$, i.e, for all $F \in$ Field, $k_{*}^{d m}(G(F)) \cong k_{*}(F)$.
km5
Theorem 4.5.5. Let $G$ be a special group. The equivalence of categories $M$ : SG $\rightarrow$ SMF induces a functor $k_{*}^{\text {mult }}: S G \rightarrow$ Igr given on the objects by $k_{*}^{\text {mult }}(G):=k_{*}^{\text {mult }}(M(G))$ and on the morphisms $f: G \rightarrow H$ by $k_{*}^{\text {mult }}(f):=k_{*}^{\text {mult }}(M(f))$. Moreover, this functor commutes with $M$ and $k^{d m}$, i.e, for all $G \in S G, k_{*}^{m u l t}(M(G)) \cong k_{*}^{d m}(G)$.

Proof. The only part requiring proof is that for all $G \in S G, k_{*}^{m u l t}(M(G)) \cong k_{*}^{d m}(G)$. The very first observation is that: since $G$ is an exponent 2 group, the reduced and non-reduced $K^{\text {mult }}$-theory of $M(G)$ coincide.

Following the argument of Theorem 2.5 in [29], it is enough to show the following two statements:
i - For all $a, b \in G$, if $b \in 1-a$ in $M(G)$ then $\lambda(b) \lambda(a)=0$;
ii - For all $a, b \in G$, if $b \in D_{G}(1, a)$ then $\rho(b) \rho(a)=0$.
For (i), if $b \in 1-a$ in $M(G)$ then $b \in D_{G}(1,-a)$ and then, $\lambda(b) \lambda(-a)=0$. Hence

$$
\lambda(b)^{2}=\lambda(b) \lambda(-a)=\lambda(b) \lambda(a)+\lambda(b) \lambda(-1) .
$$

Since $\lambda(b) \lambda(-1)=\lambda(b)^{2}$, we get $\lambda(b) \lambda(a)=0$.

For (ii) we just use the same argument: if $b \in D_{G}(1, a)$ then $b \in 1+a$ in $M(G)$ and then, $\rho(b) \rho(-a)=0$. Hence

$$
\rho(b)^{2}=\rho(b) \rho(-a)=\rho(b) \rho(a)+\rho(b) \rho(-1) .
$$

Since $\rho(b) \rho(-1)=\rho(b)^{2}$, we get $\rho(b) \rho(a)=0$.
Combining Theorems 4.5.1, 4.5.4, 4.5.5, 4.3.6 and Corollaries 4.5 .2 and 4.5 .3 we obtain the following Theorem, that unify in some sense all three K-theories:
res0
Theorem 4.5.6 (Interchanging K-theories Formulas). Let $F \in$ Field $_{2}$. Then

$$
k^{m i l}(F) \cong k^{d m}(G(F)) \cong k^{m u l t}(M(G(F))) .
$$

If $F$ is formally real and $T$ is a preordering of $F$, then

$$
k^{d m}\left(G_{T}(F)\right) \cong k^{m u l t}\left(M\left(G_{T}(F)\right)\right) .
$$

Moreover, since $M(G(F)) \cong F / m \dot{F}^{2}$ and $M\left(G_{T}(F)\right) \cong F / m T^{*}$, we get

$$
\begin{aligned}
k^{m i l}(F) & \cong k^{d m}(G(F)) \cong k^{m u l t}\left(F /{ }_{m} \dot{F}^{2}\right) \text { and } \\
k^{d m}\left(G_{T}(F)\right) & \cong k^{m u l t}\left(F /{ }_{m} T^{*}\right)
\end{aligned}
$$

Corollary 4.5.7. Let $F$ be a field. Then

$$
k^{m i l}(F) \cong k^{m u l t}\left(F /{ }_{m} \dot{F}^{2}\right)
$$

Proof. Using the previous Corollary, we already have

$$
k^{m i l}(F) \cong k^{d m}(G(F)) \cong k^{m u l t}(M(G(F))) .
$$

Now, is enough to observe that $M(G(F)) \cong F /{ }_{m} \dot{F}^{2}$.
Combining Theorem 4.5.6, Corollary 4.5.7 and Theorem 4.3.6 we get the following Corollaries.
res2
Corollary 4.5.8. Let $F$ be a formally real field and $T$ be a preordering. Then we have a surjective map

$$
k^{m i l}(F) \rightarrow k^{m u l t}\left(F /{ }_{m} T^{*}\right) .
$$

res3
Corollary 4.5.9. Let $G$ be a pre-special group and $H \subseteq G$ be a subgroup of $G$. Let $M(G)$ be the pre-special multifield associated to $G$ and $M(H)=H \cup\{0\} \subseteq M(G)$. Then

$$
G / H \cong M(G) /{ }_{m} M(H)^{*} .
$$

Moreover, $M(H) \subseteq M(G)$ is a preordering if and only if $H$ is saturated.
Corollary 4.5.10. Let $G$ be a special group and $H$ be a saturated subgroup. Then we have a surjective map

$$
k^{d m}(G) \rightarrow k^{m u l t}(G / m H) \cong k^{d m}(G / H)
$$

## Chapter 5

## Inductive Graded Rings: A Deeper Look at Marshall's Signature Conjecture

Theorem 4.5.6 gives a hint that the category of Igr is a good abstract environment for studying questions of "quadratic flavour". So a better understanding of Igr's is at least desirable and this is the main purpose of this Chapter.

We develop the general properties valid for Igr's and the main results here are Theorem 5.5.4, providing an adjunction between the categories of pre-special groups and (a subcategory of) inductive graded rings. We also characterize the Special and Weak Marshall Conjecture in the context of inductive graded rings (Section 5.6).

### 5.1 Some Categorical Facts

In order to easy the presentation, in this section there are some categorical results concerning adjunctions. Mostly are based on [8], but the reader could also consult [44].
3.1.1borceux

Definition 5.1.1 (3.1.1 of [8]). Let $F: \mathcal{A} \rightarrow \mathcal{B}$ be a functor and $B$ an object of $\mathcal{B}$. A reflection of $B$ along $F$ is a pair $\left(R_{B}, \eta_{B}\right)$ where

1. $R_{B}$ is an object of $A$ and $\eta_{B}: B \rightarrow F\left(R_{B}\right)$ is a morphism of $\mathcal{B}$.
2. If $A \in \mathcal{A}$ is another object and $b: B \rightarrow F(A)$ is a morphism of $\mathcal{B}$, there exists a unique morphism $a: R_{B} \rightarrow A$ in $\mathcal{A}$ such that $F(a) \circ \eta_{B}=b$.
3.1.2borceux

Proposition 5.1.2 (3.1.2 of [8]). Let $F: \mathcal{A} \rightarrow \mathcal{B}$ be a functor and $B$ an object of $\mathcal{B}$. When the reflection of $B$ along $F$ exists, it is unique up to isomorphism.
3.1.4borceux

Definition 5.1.3 (3.1.4 of [8]). A functor $R: \mathcal{B} \rightarrow A$ is left adjoint to the functor $F: \mathcal{A} \rightarrow \mathcal{B}$ when there exists a natural transformation

$$
\eta: 1_{\mathcal{B}} \Rightarrow F \circ R
$$

such that for every $B \in \mathcal{B},\left(R(B), \eta_{B}\right)$ is a reflection of $B$ along $F$.
Theorem 5.1.4 (3.1.5 of [8]). Consider two functors $F: \mathcal{A} \rightarrow \mathcal{B}$ and $G: \mathcal{B} \rightarrow \mathcal{A}$. The following conditions are equivalent.

1. $G$ is left adjoint of $F$.
2. There exist a natural transformation $\eta: 1_{\mathcal{B}} \Rightarrow F \circ G$ and $\varepsilon: G \rightarrow F \Rightarrow 1_{\mathcal{A}}$ such that

$$
(F * \varepsilon) \circ(\eta * F)=1_{F},(\varepsilon * G) \circ(G * \eta)=1_{G} .
$$

3. There exist bijections

$$
\theta_{A B}: \mathcal{A}(G(B), A) \cong \mathcal{B}(B, F(A))
$$

for every objects $A$ and $B$, and those bijections are natural both in $A$ and $B$.
4. $F$ is right adjoint of $G$.
3.2.2borceux

Proposition 5.1.5 (3.2.2 of [8]). If the functor $F: \mathcal{A} \rightarrow \mathcal{B}$ has a left adjoint then $F$ preserves all limits which turn out to exist in $\mathcal{A}$.
3.4.1borceux

Proposition 5.1.6 (3.4.1 of [8]). Consider two functors $F: \mathcal{A} \rightarrow \mathcal{B}, G: \mathcal{B} \rightarrow \mathcal{A}$ with $G$ left adjoint to $F$ with $\eta: 1_{\mathcal{B}} \Rightarrow F \circ G$ and $\varepsilon: G \circ F \Rightarrow 1_{\mathcal{A}}$ the two corresponding natural transformations. The following conditions are equivalent.

1. F is full and faithfull.
2. $\varepsilon$ is an isomorphism.

Under these conditions, $\eta * F$ and $G * \eta$ are isomorphisms as well.
3.4.3borceux

Proposition 5.1.7 (3.4.3 of [8]). Given a functor $F: \mathcal{A} \rightarrow \mathcal{B}$, the following conditions are equivalent:

1. $F$ is full and faithfull and has a full and faithfull left adjoint $G$.
2. F has a left adjoint $G$ and the two canonical natural transformations of the adjunction $\eta$ : $1_{\mathcal{B}} \Rightarrow F \circ G$ and $\varepsilon: G \rightarrow F \Rightarrow 1_{\mathcal{A}}$ are isomorphisms.
3. There exists a functor $G: \mathcal{B} \rightarrow \mathcal{A}$ and two arbitrary natural isomorphisms $1_{\mathcal{B}} \cong F \circ G$, $G \circ F \cong 1_{\mathcal{A}}$.
4. $F$ is full and faithfull and each object $B \in \mathcal{B}$ is isomorphic to an object of the form $F(A)$, for some $A \in \mathcal{A}$.
5. The dual condition of (1).
6. The dual condition of (2).
3.4.4borceux

Definition 5.1.8 (3.4.4 of [8]). The categories $\mathcal{A}, \mathcal{B}$ are equivalent if there exist a functor $F$ : $\mathcal{A} \rightarrow \mathcal{B}$ satisfying the conditions of Proposition 5.1.7.

### 5.2 The First Properties of Igr

In this section we discuss the theory of Igr's. Constructions like products, limits, colimits, ideals, quotients, kernel and image are not new and are obtained in a very straightforward manner (basically, putting those structures available for rings in a "coordinatewise" fashion), then in order to gain speed, we will present these facts leaving more detailed proofs to the reader.

Denote: $p \mathbb{F}_{2}-\bmod$ the category of pointed $\mathbb{F}_{2}$-modules, Ring the category of commutative rings with unity and morphism that preserves these units and $\mathrm{Ring}_{2}$ the full subcategory of the associative $\mathbb{F}_{2}$-algebras. We have a functorial correspondence $\operatorname{Ring}_{2} \rightarrow \mathrm{Igr}$, given by the following diagram:


Here $A$ is a $p \mathbb{F}_{2}-\bmod$ where $\top_{n}=1, n \geq 1$ and $\top_{0}=1 \in \mathbb{F}_{2}$.
trivialigr
Definition 5.2.1. The trivial graded ring functor $\mathbb{T}:$ Ring $_{2} \rightarrow$ Igr is the functor defined for $f: A \rightarrow B$ by $T(A)_{0}:=\mathbb{F}_{2}, T(f)_{0}:=i d_{\mathbb{F}_{2}}$ and for all $n \geq 1$ we set $T(A)_{n}=A$ and $T(f)_{n}:=f$.
f2alg
Definition 5.2.2. We define the associated $\mathbb{F}_{2}$-algebra functor $\mathbb{A}:$ Igr $\rightarrow$ Ring $g_{2}$ is the functor defined for $f: R \rightarrow S$ by

$$
\mathbb{A}(R):=R_{\mathbb{A}}=\underset{n \geq 0}{\lim } R_{n} \text { and } \mathbb{A}(f)=f_{\mathbb{A}}:=\underset{n \geq 0}{\lim _{n}} f_{n} .
$$

More explicitly, $\mathbb{A}(R)=\left(R_{\mathbb{A}}, 0,1,+_{\mathbb{A}}, \cdot\right)$, where
$\mathrm{i}-R_{\mathbb{A}}=\underset{n \geq 0}{\lim } R_{n}$,
ii - $0=[(0,0)]$ and $1=[(1,0)]$,
iii - given $\left[\left(a_{n}, n\right)\right],\left[\left(b_{m}, m\right)\right] \in R_{\mathbb{A}}$ and setting $d \geq m, n$ we have

$$
\left[\left(a_{n}, n\right)\right]+\left[\left(b_{m}, m\right)\right]=\left[\left(h_{n d}\left(a_{n}\right)+h_{m d}\left(b_{m}\right), d\right)\right]
$$

iv - given $\left[\left(a_{n}, n\right)\right],\left[\left(b_{m}, m\right)\right] \in R_{\mathbb{A}}$, we have

$$
\left[\left(a_{n}, n\right)\right] \cdot\left[\left(b_{m}, m\right)\right]=\left[\left(a_{n} *_{n m} b_{m}, n+m\right)\right]
$$

propadj1

## Proposition 5.2.3.

$i$ - The functor $\mathbb{A}$ is the left adjunct to $\mathbb{T}$.
ii - The functor $\mathbb{T}$ is full and faithful.
iii - The composite functor $\mathbb{A} \circ \mathbb{T}$ is naturally isomorphic to the functor $1_{R i n g}{ }_{2}$.
Proof. Let $R \in I g r$. We have

$$
\mathbb{T}(\mathbb{A}(R))=\mathbb{T}\left(\underset{m \geq 0}{\lim } R_{m}\right)
$$

In other words, for all $n \geq 1$

$$
\mathbb{T}\left(\underset{m \geq 0}{\lim } R_{m}\right)_{n}:=\underset{m \geq 0}{\lim _{m}} R_{m}
$$

Then, for all $n \geq 1$ we have a canonical embedding

$$
\eta(R)_{n}: R_{n} \rightarrow{\underset{m \geq 0}{ }}_{\lim _{m \geq 0}} R_{m}=\mathbb{T}\left({\underset{m \geq 0}{\lim }}_{{ }_{m \geq}} R_{m}\right)_{n}
$$

providing a morphism

$$
\eta(R): R \rightarrow \underset{m \geq 0}{\lim } R_{m}=\mathbb{T}\left(\underset{m \geq 0}{\lim } R_{m}\right) .
$$

For $f \in \operatorname{Igr}(R, S)$, taking $n \geq 1$ we have a commutative diagram

with the convention that $\eta(R)_{0}=i d_{\mathbb{F}_{2}}$. Then it is legitimate the definition of a natural transformation $\eta: 1_{\mathrm{Igr}} \rightarrow \mathbb{T} \circ \mathbb{A}$ given by the rule $R \mapsto \eta(R)$.

Now let $A \in \operatorname{Ring}_{2}$ and $g \in \operatorname{Ring}_{2}(R, \mathbb{T}(A))$. Then for each $n \geq 0$, there is a morphism $g_{n}: R_{n} \rightarrow \mathbb{T}(A)_{n}=A$ and by the universal property of inductive limit we get a morphism

$$
\underset{m \geq 0}{\lim _{\geq}} g_{n}: \varliminf_{m \geq 0} R_{m} \rightarrow A .
$$

In fact, $\underset{m \geq 0}{\lim _{\longrightarrow 0}} g_{n}=\mathbb{A}(g)$.

Now, using the fact that $\eta(R)_{n}$ is the morphism induced by the inductive limit we have for all $n \geq 0$ the following commutative diagram


In other words, $\eta(B)_{n}$ is the canonical morphism commuting the diagram

and hence, $\mathbb{A}$ is the left adjoint of $\mathbb{T}$, proving item (i). By the very definition of $\mathbb{A}$ and $\mathbb{T}$ we get item (iii), and using Proposition 5.1.5 we get item (ii).

Using Proposition 5.1.5 (and its dual version) we get the following Corollary.

## Corollary 5.2.4.

$i-\mathbb{T}:$ Ring $_{2} \rightarrow$ Igr preserves all projective limits.
ii - If I is such that Igr is I-inductively complete then for $\left\{A_{i}\right\}_{i \in I}$ in Igr we have

$$
\lim _{i \in I} A_{i} \cong \mathbb{A}\left(\lim _{i \in I} \mathbb{T}\left(A_{i}\right)\right) .
$$

iii $-\mathbb{F}_{2} \in$ Ring $_{2}$ is the initial object in Ring $_{2}$.
iv - $0 \in \operatorname{Ring}_{2}$ is the terminal object in $\operatorname{Ring}_{2}$.
$v-\mathbb{T}\left(\mathbb{F}_{2}\right)$ is the initial object in Igr.
$v i-\mathbb{T}(0)$ is the terminal object in Igr.
Now we discuss (essentially) the limits and colimits in Igr. Fix a non-empty set $I$ and let $\left\{\left(R_{i}, \top_{i}, h_{i}\right)\right\}_{i \in I}$ be a family of Igr's. We start with the construction of the Igr-product

$$
R=\prod_{i \in I} R_{i} .
$$

For this, we define $R_{0} \cong \mathbb{F}_{2}$ and for all $n \geq 1$, we define

$$
R_{n}:=\prod_{i \in I}\left(R_{i}\right)_{n} \text { and } \top_{n}:=\prod_{i \in I}\left(\top_{i}\right)_{n} .
$$

In the sequel, we define $h_{0}: \mathbb{F}_{2} \rightarrow R_{1}$ as the only possible morphism and for $n \geq 1$, we define $h_{n}: R_{n} \rightarrow R_{n+1}$ by

$$
h_{n}:=\prod_{i \in I}\left(h_{i}\right)_{n}
$$

## Definition 5.2.5.

$i$ - The space of orderings, $X_{R}$, of the Igr $R$, is the set of Igr-morphisms $\operatorname{Igr}\left(R, \mathbb{T}\left(\mathbb{F}_{2}\right)\right.$. By the Proposition 5.2.3. ( $i$ ), we have a natural bijection $\operatorname{Igr}\left(R, \mathbb{T}\left(\mathbb{F}_{2}\right) \cong \operatorname{Ring}_{2}\left(\mathbb{A}(R), \mathbb{F}_{2}\right)\right.$, thus considering the discrete topologies on the $\mathbb{F}_{2}$-algebras $\left.\mathbb{A}(R), \mathbb{F}_{2}\right)$ and transporting the boolean topology in $\operatorname{Ring}_{2}\left(\mathbb{A}(R), \mathbb{F}_{2}\right)$, we obtain a boolean topology on the space of orderings $X_{R}=$ $\operatorname{Igr}\left(R, \mathbb{T}\left(\mathbb{F}_{2}\right)\right)$.
ii- The boolean hull, $B(R)$, of the Igr $R$, is the boolean ring canonically associated to the space of orderings of $R$ by Stone duality: $B(R):=\mathcal{C}\left(X_{R}, \mathbb{F}_{2}\right)$.
iii- $A$ Igr $R$ is called formally real if $X_{R} \neq \emptyset$ (or, equivalently, if $B(R) \neq 0$ ).
fixigr1
Proposition 5.2.6. Let $I$ be a non-empty set and $\left\{\left(R_{i}, h_{i}\right)\right\}_{i \in I}$ be a family of Igr's. Then

$$
R=\prod_{i \in I} R_{i}
$$

with the above rules is an Igr. Moreover it is the product in the category Igr.
Proof. Using Definition 4.4.1 is straightforward to verify that $\left(R, \top_{n}, h_{n}\right)$ is an Igr. Note that for each $i \in I$, we have an epimorphism $\pi_{i}: R \rightarrow R_{i}$ given by the following rules: for each $n \geq 0$ and each $\left(x_{i}\right)_{i \in I} \in R_{n}$, we define

$$
\left(\pi_{i}\right)_{n}\left(\left(x_{i}\right)_{i \in I}\right):=x_{i} .
$$

Now, let $\left(Q,\left\{q_{i}\right\}_{i \in I}\right)$ be another pair with $Q$ being an Igr and $q_{i}: Q \rightarrow R_{i}$ being a morphism for each $i \in I$. Given $i \in I$ and $n \geq 0$, since $R_{n}:=\prod_{i \in I}\left(R_{i}\right)_{n}$ is the product in the category of pointed $\mathbb{F}_{2}$-modules, we have an unique morphism $(q)_{n}:(Q)_{n} \rightarrow(R)_{n}$ such that $\left(\pi_{i}\right)_{n} \circ(q)_{n}=\left(q_{i}\right)_{n}$. Set $q_{n}:=\left(\left(q_{i}\right)_{i \in I}\right)_{n}$. By construction, $q$ is the unique Igr-morphism such that $\pi_{i} \circ q=q_{i}$, completing the proof that $R$ is in fact the product in the category Igr.

## Proposition 5.2.7.

$i$ - Let $R$ be an Igr and let $X \subseteq R=\bigoplus_{n \in \mathbb{N}} R_{n}$. Then there exists the inductive graded subring generated by $X$ (notation : $[X] \stackrel{i}{\hookrightarrow} R$ ): this is the least inductive graded subring of $R$ such that $\forall n \in \mathbb{N}, X \cap R_{n} \subseteq[X]_{n}$.
ii- Let $\mathcal{I}$ be a small category and $\mathcal{R}: \mathcal{R} \rightarrow$ Igr be a diagram. Then there exists $\lim _{\underset{i}{ } \in \mathcal{I}} \mathcal{R}_{i}$ in the category Igr.

Proof.
i- It is enough consider $S_{X}$, the $\mathbb{F}_{2}$-subalgebra of $\left(\bigoplus_{n \in \mathbb{N}} R_{n}, *\right)$ generated by $X \cup\left\{T_{1}\right\} \subseteq$ $\bigoplus_{n \in \mathbb{N}} R_{n}$ and set $\forall n \in \mathbb{N},[X]_{n}:=s_{x} \cap R_{n}$.
ii- Just define $\varliminf_{i \in \mathcal{I}} \mathcal{R}_{i}$ as the inductive graded subring of $\prod_{i \in o b j(\mathcal{I})} \mathcal{R}_{i}$ generated by $X_{D}=$ $\bigoplus_{n \in \mathbb{N}} X_{n}$ and $X_{n}:=\varliminf_{i \in \mathcal{I}}\left(\mathcal{R}_{i}\right)_{n}$ (projective limit of pointed $\mathbb{F}_{2}$-algebras).

Now we construct the Igr-tensor product of a finite family of Igr's, $\left\{R_{i}: i \in I\right\}$

$$
R=\bigotimes_{i \in I} R_{i} .
$$

For this, we define $R_{0} \cong \mathbb{F}_{2}$ and for all $n \geq 1$, we define

$$
\begin{gathered}
R_{n}:=\bigotimes_{i \in I}\left(R_{i}\right)_{n}, \\
\left(\otimes_{i \in I} a_{i}\right) *_{n, k}\left(\otimes_{i \in I} b_{i}\right):=\otimes_{i \in I}\left(a_{i} *_{n, k}^{i} b_{i}\right) \\
\text { and } \top_{n}:=\otimes_{i \in I}\left(\top_{i}\right)_{n} .
\end{gathered}
$$

In particular, if $I=\emptyset$, then $R_{n}=\{0\}, n \geq 1$. In the sequel, we define $h_{0}: \mathbb{F}_{2} \rightarrow R_{1}$ as the only possible morphism and for $n \geq 1$, we define $h_{n}: R_{n} \rightarrow R_{n+1}$ by

$$
h_{n}:=\bigotimes_{i \in I}\left(h_{i}\right)_{n}
$$

In other words, for a generator $\bigotimes_{i \in I} x_{i} \in R_{n}$, we have

$$
h_{n}\left(\otimes_{i \in I} x_{i}\right):=\bigotimes_{i \in I}\left(h_{i}\right)_{n}\left(x_{i}\right)
$$

Proposition 5.2.8. Let $I$ be a finite set and $\left\{\left(R_{i}, h_{i}\right)\right\}_{i \in I}$ be a family of Igr's. Then

$$
R=\bigotimes_{i \in I} R_{i}
$$

with the above rules is an Igr. Moreover it is the coproduct in the category Igr.
Now suppose that $(I, \leq)$ is an upward directed poset and that $\left(\left(R_{i}, h_{i}\right), \varphi_{i j}\right)_{i \leq j \in I}$ is an inductive system of Igr's. We define the inductive limit

$$
R=\underset{i \in I}{\lim } R_{i}
$$

by the following: for all $n \geq 0$ define

$$
R_{n}:=\underset{i \in I}{\lim }\left(R_{i}\right)_{n} .
$$

Note that

$$
R_{0}:=\underset{i \in I}{\lim }\left(R_{i}\right)_{0} \cong \underset{i \in I}{\lim } \mathbb{F}_{2} \cong \mathbb{F}_{2} .
$$

In the sequel, for $n \geq 1$ we define $h_{n}: R_{n} \rightarrow R_{n+1}$ by

$$
h_{n}:=\underset{i \in I}{\lim }\left(h_{i}\right)_{n}
$$

Proposition 5.2.9. Let $(I, \leq)$ is an upward directed poset and $\left(\left(R_{i}, h_{i}\right), \varphi_{i j}\right)_{i \in I}$ be a directed family of Igr's. Then

$$
R=\underset{i \in I}{\lim } R_{i}
$$

with the above rules is an Igr. Moreover, it is the inductive limit in the category Igr.
Proposition 5.2.10. The general coproduct (general tensor product) of a family $\left\{R_{i}: i \in I\right\}$ in
the category Igr is given by the combination of constructions:

$$
\bigotimes_{i \in I} R_{i}:=\underset{I^{\prime} \in \underset{P_{f i n}}{ }(I)}{\lim _{i \in I^{\prime}}} \bigotimes_{i}
$$

After discussing directed inductive colimits and coproducts, we will deal with ideals, quotients, and coequalizers.

Definition 5.2.11. Given $R \in \operatorname{Igr}$ and $\left(J_{n}\right)_{n \geq 0}$ where $J_{n} \subseteq R_{n}$ for all $n \geq 0$. We say that $J$ is a graded ideal of $R$ where

$$
J:=\bigoplus_{n \geq 0} J_{n} \subseteq \bigoplus_{n \geq 0} R_{n}
$$

is an ideal of $(R, *)$.
In particular, for all $n \geq 0, J_{n} \subseteq R_{n}$ is a graded $\mathbb{F}_{2}$-submodule of $\left(R_{n},+_{n}, 0_{n}\right)$. For each $X \subseteq R$, there exists the ideal generated by $X$, denoted by $\langle X\rangle$. It is the smaller graded ideal of $R$ such that for all $n \geq 0,\left(X \cap R_{n}\right) \subseteq[X]_{n}$. For this, just consider $\langle X\rangle$, the ideal of $(R, *)$ generated by $X \subseteq R$ and define $\langle X\rangle_{n}:=\langle X\rangle \cap R_{n}$.

Definition 5.2.12. Let $R, S$ be Igr's and $f: R \rightarrow S$ be a morphism. We define the kernel of $f$, notation $\operatorname{Ker}(f)$ by

$$
\operatorname{Ker}(f)_{n}:=\left\{x \in R_{n}: f_{n}(x)=0\right\}
$$

and image of $f$, notation $\operatorname{Im}(f)$ by

$$
\operatorname{Im}(f)_{n}:=\left\{f_{n}(x) \in S_{n}: x \in R_{n}\right\}
$$

Of course, $\operatorname{Ker}(f) \subseteq R$ is an ideal and $\operatorname{Im}(f) \subseteq S$ is an Igr.
Given $R \in \operatorname{Igr}$ and $J=\left(J_{n}\right)_{n \geq 0}$ a graded ideal of $R$, we define $R / J \in \operatorname{Igr}$, the quotient inductive graded ring of $R$ by $J$ : for all $n \geq 0,(R / J)_{n}:=R_{n} / J_{n}$, where the distinguished element is $\top_{n}+_{n} J_{n}$. We have a canonical projection $q_{J}: R \rightarrow R / J$, "coordinatewise surjective" and therefore, an Igr-epimorphism.

Proposition 5.2.13 (Homomorphism Theorem). Let $R, S$ be Igr's and $f: R \rightarrow S$ be a morphism. Then there exist an unique monomorphism $\bar{f}: R / \operatorname{Ker}(f) \rightarrow S$ commuting the following diagram:

where $q$ is the canonical projection. In particular $R / \operatorname{Ker}(f) \cong \operatorname{Im}(f)$.
Proposition 5.2.14. Let $R \underset{g}{\stackrel{f}{\rightrightarrows}} S$ be Igr-morphisms and consider $q_{J}: S \rightarrow S / J$ the quotient morphism where $J:=\langle X\rangle$ is the graded ideal generated by $X_{n}:=\left\{f_{n}(a)-g_{n}(a): a \in R_{n}\right\}, n z i n \mathbb{N}$. Then $q_{J}$ is the coequalizer of $f, g$.

Proposition 5.2.15. Given $R, S \in \operatorname{Igr}$ and $f \in \operatorname{Igr}(R, S)$.
$i$ - $f$ is a Igr-monomorphism whenever for all $n \geq 0 f_{n}: R_{n} \rightarrow S_{n}$ is a monomorphism of pointed $\mathbb{F}_{2}$-modules iff for all $n \geq 0, f_{n}: R_{n} \rightarrow S_{n}$ is an injective homomorphism of pointed $\mathbb{F}_{2}$-modules.
ii - $f$ is a Igr-epimorphism whenever for all $n \geq 0 f_{n}: R_{n} \rightarrow S_{n}$ is a epimorphism of pointed $\mathbb{F}_{2}$ modules iff for all $n \geq 0, f_{n}: R_{n} \rightarrow S_{n}$ is a surjective homomorphism of pointed $\mathbb{F}_{2}$-modules.
iii - $f$ is a Igr-isomorphism iff for all $n \geq 0 f_{n}: R_{n} \rightarrow S_{n}$ is a isomorphism of pointed $\mathbb{F}_{2}$-modules iff for all $n \geq 0, f_{n}: R_{n} \rightarrow S_{n}$ is a bijective homomorphism of pointed $\mathbb{F}_{2}$-modules.
Definition 5.2.16. We denote Igr $_{\text {fin }}$ the full subcategory of Igr such that

$$
\operatorname{Obj}\left(\operatorname{Igr}_{f i n}\right)=\left\{R \in \operatorname{Obj}(\operatorname{Igr}):\left|R_{n}\right|<\omega \text { for all } n \geq 1\right\} .
$$

Remark 5.2.17. Of course,

$$
\left\{R \in \operatorname{Obj}(I g r):\left|\bigoplus_{n \geq 1} R_{n}\right|<\omega\right\} \neq \operatorname{Obj}\left(\operatorname{Igr}_{f i n}\right),
$$

for example, in 4.4.2 (a), if $F$ is a Euclidian field (for instance, any real closed field), then $\bigoplus_{n \in \mathbb{N}} I^{n} F / I^{n+1} F$
 finite.

### 5.3 Relevant subcategories of Igr

The aim of this Section is to define subcategories of Igr that mimetize the following two central aspects of K-theories:

1. The K-theory graded ring is "generated" by $K_{1}$;
2. The K-theory graded ring is defined by some convenient quotient of a graded tensor algebra.

Our desired category will be the intersection of two subcategories. The first one is obtained after we define the graded subring generated by the level 1 functor

$$
\mathbb{1}: \operatorname{Igr} \rightarrow \mathrm{Igr} .
$$

We define it as follow: for an object $R=\left(\left(R_{n}\right)_{n \geq 0},\left(h_{n}\right)_{n \geq 0}, *_{n m}\right)$,
i - $\mathbb{1}(R)_{0}:=R_{0} \cong \mathbb{F}_{2}$,
ii - $\mathbb{1}(R)_{1}:=R_{1}$,
iii - for $n \geq 2$,

$$
\begin{aligned}
\mathbb{1}(R)_{n} & :=\left\{x \in R_{n}: x=\sum_{j=1}^{r} a_{1 j} *_{11} \ldots *_{11} a_{n j}\right. \\
& \text { with } \left.a_{i j} \in R_{1}, 1 \leq i \leq n, 1 \leq j \leq r \text { for some } r \geq 1\right\} .
\end{aligned}
$$

Note that for all $n \geq 2, R_{n}$ is generated by the expressions of type

$$
d_{1} *_{11} d_{2} *_{11} \ldots *_{11} d_{n}, d_{i} \in R_{1}, i=1, \ldots, n
$$

Of course, $\mathbb{1}(R)$ provides an inclusion $\iota_{\mathbb{1}(R)}: \mathbb{1}(R) \rightarrow R$ in the obvious way.
On the morphisms, for $f \in \operatorname{Igr}(R, S)$, we define $\mathbb{1}(f) \in \operatorname{Igr}(\mathbb{1}(R), \mathbb{1}(S))$ by the restriction $\mathbb{1}(f)=f 1_{\mathbb{1}(R)}$. In other words, $\mathbb{1}(f)$ is the only Igr-morphisms that makes the following diagram commute:


Definition 5.3.1. We denote $I g r_{\mathbb{1}}$ the full subcategory of Igr such that

$$
\operatorname{Obj}\left(\operatorname{Igr_{\mathbb {1}}}\right)=\left\{R \in \operatorname{Igr}: \iota_{\mathbb{1}(R)}: \mathbb{1}(R) \rightarrow R \text { is an isomorphism }\right\} .
$$

## Example 5.3.2.

$i$ - If $A$ is a $\mathbb{F}_{2}$-algebra, then $\mathbb{T}(A) \in \operatorname{obj}\left(\operatorname{Igr_{\mathbb {\Perp }}}\right)$.
ii - If $F$ is an hyperbolic hyperfield, then $k_{*}(F) \in \operatorname{obj}\left(\operatorname{Igr}_{\mathbb{1}}\right)$.
iii - If $F$ is a special hyperfield (equivalently, $G=F \backslash\{0\}$ is a special group), then the graduate Witt ring of $F$ (definition 5.4.9) $W_{*}(F) \in \operatorname{obj}\left(\operatorname{Igr}_{\mathbb{1}}\right)$.
iv - If $F$ is a field with char $(F) \neq 2$, then, by a known result of Vladimir Voevodski,

$$
\mathcal{H}^{*}\left(\operatorname{Gal}\left(F^{s} \mid F\right),\{ \pm 1\}\right) \in \operatorname{obj}\left(\operatorname{Igr_{\mathbb {1}}}\right) .
$$

## Proposition 5.3.3.

$i$ - For each $R \in$ Igr we have that $\iota_{\mathbb{1}(\mathbb{1}(R))}: \mathbb{1}(\mathbb{1}(R)) \rightarrow \mathbb{1}(R)$ is the identity arrow.
$i i-\mathbb{1} \circ \mathbb{1}=\mathbb{1}$.
iii - The functor $\mathbb{1}: I g r \rightarrow I g r_{\mathbb{1}}$ is the right adjoint of the inclusion functor $j_{\mathbb{1}}: I g r_{\mathbb{1}} \rightarrow I g r$.
$i v-j_{\mathbb{1}}: I g r_{\mathbb{1}} \rightarrow$ Igr creates inductive limits and to obtain the projective limits in Igr $r_{\mathbb{1}}$ is sufficient restrict the projective limits obtained in Igr:

$$
\lim _{\overparen{i} \in I} R_{i} \cong\left(\lim _{\underset{i \in I}{ }} j_{\mathbb{I}}\left(R_{i}\right)\right)_{\mathbb{1}} \xrightarrow{\lim _{i \in I} j_{\mathbb{1}}\left(R_{i}\right)} \lim _{\overparen{i \in I}} j_{\mathbb{B}}\left(R_{i}\right) .
$$

Proof. Similar to Proposition 5.2.3.
Now we define the second subcategory. We define the quotient graded ring functor

$$
\mathcal{Q}: \operatorname{Igr} \rightarrow \mathrm{Igr}
$$

as follow: for a object $R=\left(\left(R_{n}\right)_{n \geq 0},\left(h_{n}\right)_{n \geq 0}, *_{n m}\right), \mathcal{Q}(R):=R / T$, where $T=\left(T_{n}\right)_{n \geq 0}$ is the ideal generated by $\left\{\left(T_{1}+{ }_{1} a\right) *_{11} a \in R_{2}: a \in R_{1}\right\}$. More explicit,
i - $T_{0}:=\left\{0_{0}\right\} \subseteq R_{0}$,
ii - $T_{1}:=\left\{0_{1}\right\} \subseteq R_{1}$,
iii - for $n \geq 2, T_{n} \subseteq R_{n}$ is the pointed $\mathbb{F}_{2}$-submodule generated by

$$
\begin{aligned}
& \left\{x \in R_{n}: x=y_{l} *_{l 1}\left(\top_{1}+{ }_{1} a_{1}\right) *_{11} a_{1} *_{1 r} z_{r},\right. \\
& \left.\quad \text { with } a_{1} \in R_{1}, y_{l} \in R_{l}, z_{r} \in R_{r}, l+r=n-2\right\} .
\end{aligned}
$$

Of course, $\mathcal{Q}(R)$ provides a projection $\pi_{R}: R \rightarrow \mathcal{Q}(R)$ in the obvious way.
On the morphisms, for $f \in \operatorname{Igr}(R, S)$, we define $\mathcal{Q}(f) \in \operatorname{Igr}(\mathcal{Q}(R), \mathcal{Q}(S))$ by the only Igrmorphisms that makes the following diagram commute:

quotop
Definition 5.3.4. We denote Igr $_{h}$ the full subcategory of Igr such that

$$
\operatorname{Obj}\left(\operatorname{Igr} r_{h}\right)=\left\{R \in \operatorname{Igr}: \pi_{R}: R \rightarrow \mathcal{Q}(R) \text { is an isomorphism }\right\} .
$$

Remark 5.3.5. Note that $R \in \operatorname{obj}\left(\right.$ Igr $\left._{h}\right)$ iff for each $a \in R_{1}, a *_{11} \top_{1}=a *_{11} a \in R_{2}$. Each $R$ satisfying this condition is, in some sense, "hyperbolic" (see Proposition 5.5.2): this is the motivation of the index " $h$ ".
Example 5.3.6. $\quad i$ - Let $A$ be a $\mathbb{F}_{2}$-algebra. Then $\mathbb{T}(A) \in o b j\left(I g r_{h}\right)$ iff $A$ is a boolean ring (i.e., $\left.\forall a \in A, a^{2}=a\right)$.
ii- If $F$ is an hyperbolic hyperfield, then $k_{*}(F) \in \operatorname{obj}\left(I g r_{h}\right)$.
iii- If $F$ is a special hyperfield (equivalently, $G=F \backslash\{0\}$ is a special group), then $W_{*}(F) \in$ obj $\left(I g r_{h}\right)$.
iv- If $F$ is a field with char $(F) \neq 2$, then $\mathcal{H}^{*}\left(\operatorname{Gal}\left(F^{s} \mid F\right),\{ \pm 1\}\right) \in \operatorname{obj}\left(\operatorname{Igr}_{h}\right)$.

## Proposition 5.3.7.

$i$ - For each $R \in$ Igr we have that $\pi_{\mathcal{Q}(R)}: \mathcal{Q}(R) \rightarrow \mathcal{Q}(\mathcal{Q}(R))$ is an isomorphism.
$i i-\mathcal{Q} \circ \mathcal{Q}=\mathcal{Q}$.
iii - The functor $\mathcal{Q}: I g r \rightarrow I g r_{h}$ is the left adjoint of the inclusion functor $j_{q}: I g r_{\mathcal{Q}} \rightarrow I g r$.
iv - $j_{q}: I g r_{h} \rightarrow$ Igr creates projective limits and to obtain the inductive limits in Igr ${ }_{h}$ is sufficient restrict the inductive limits obtained in Igr:

$$
\lim _{i \in I} j_{q}\left(R_{i}\right) \xrightarrow{\underset{\overrightarrow{i \in I}}{\lim _{q} j_{q}\left(R_{i}\right)}}\left(\underset{\vec{Q} \in I}{ } j_{q}\left(R_{i}\right)\right)_{\mathcal{Q}}^{\cong} \lim _{i \in I} R_{i} .
$$

Moreover, $j_{q}: I g r_{h} \rightarrow$ Igr creates filtered inductive limits and quotients by graded ideals.
Are examples of inductive graded rings in $I g r_{+}$: (i) $\mathbb{T}(A)$, where $A$ is a boolean ring; (ii) $k_{*}(F)$, where $F$ is an hyperbolic hyperfield; (iii) $W_{*}(F)$, where $F$ is an special hyperfield; (iv) $\mathcal{H}^{*}\left(\operatorname{Gal}\left(F^{s} \mid F\right),\{ \pm 1\}\right)$, where $F$ is a field with $\operatorname{char}(F) \neq 2$.
igr+
Definition 5.3.8 (The Category $\mathrm{Igr}_{+}$). We denote by $\mathrm{Igr}_{+}$the full subcategory of Igr such that

$$
O b j\left(I g r_{+}\right)=\operatorname{Obj}\left(I g r_{\mathbb{1}}\right) \cap \operatorname{Obj}\left(I g r_{h}\right) .
$$

We denote by $j_{+}:$Igr $_{+} \rightarrow$ Igr the inclusion functor.
Remark 5.3.9. $i$ - Note that the notion of an Igr, R, be in the subcategory Igr ${ }_{h}$ can be axiomatized by a first-order (finitary) sentence in $L$, the polysorted language for Igr's described in the previous Chapter: $\left(\forall a: 1, a *_{11} a=\top_{1} *_{11} a\right)$. On the other hand, the concepts $R \in I^{\prime} r_{\mathbb{1}}$ and $R \in I g r_{+}$are axiomatized by $L_{\omega_{1}, \omega}$-sentences.
ii- Note that the subcategory $\mathrm{Igr}_{+} \hookrightarrow$ Igr is closed by filtered inductive limits.
In order to think of an object in $\mathrm{Igr}_{+}$as a graded ring of "K-theoretic type", we make the following convention.
igrlog
Definition 5.3.10 (Exponential and Logarithm of an $\operatorname{Igr}$ ). Let $R \in I g r_{+}$and write $R_{1}$ multiplicatively by $(\Gamma(R), \cdot, 1,-1)$, i.e, fix an isomorphism $e_{R}: R_{1} \rightarrow \Gamma(R)$ in order that $e_{R}(T)=-1$ and $e_{R}(a+b)=a \cdot b$. Such isomorphism $e_{R}$ is called exponential of $R$ and $l_{R}=e_{R}^{-1}$ is called logarithm of $R$. In this sense, we can write $R_{1}=\{l(a): a \in \Gamma(R)\}$. We also denote $l(a) *_{11} l(b)$ simply by $l(a) l(b), a, b \in \Gamma(R)$. We drop the superscript and write just $e, l$ when the context allows $i t$.

Using Definitions 5.3.8, 5.3.10 (and of course, Definitions 5.3.1 and 5.3.4 with an argument similar to the used in Lemma 4.3.2) we have the following properties.

Lemma 5.3.11 (First Properties). Let $R \in I g r_{+}$.

$$
i-l(1)=0 .
$$

ii - For all $n \geq 1, \eta \in R_{n}$ is generated by $l\left(a_{1}\right) \ldots l\left(a_{n}\right)$ with $a_{1}, \ldots, a_{n} \in \Gamma(R)$.
iii $-l(a) l(-a)=0$ and $l(a) l(a)=l(-1) l(a)$ for all $a \in \Gamma(R)$.
$i v-l(a) l(b)=l(b) l(a)$ for all $a, b \in \Gamma(R)$.
$v$ - For every $a_{1}, \ldots, a_{n} \in \Gamma(R)$ and every permutation $\sigma \in S_{n}$,

$$
l\left(a_{1}\right) \ldots l\left(a_{i}\right) \ldots l\left(a_{n}\right)=\operatorname{sgn}(\sigma) l\left(a_{\sigma 1}\right) \ldots l\left(a_{\sigma n}\right) \text { in } R_{n} .
$$

vi- For all $\xi \in R_{n}, \eta \in R_{n}$,

$$
\xi \eta=\eta \xi
$$

vii- For all $n \geq 1$,

$$
h_{n}\left(l\left(a_{1}\right) \ldots l\left(a_{n}\right)\right)=l(-1) l\left(a_{1}\right) \ldots l\left(a_{n}\right) .
$$

Proposition 5.3.12. Let $R \in$ Igr $_{+}$
$i$ - For each $n \in \mathbb{N}$ and each $x \in R_{n}, x *_{n, n} x=\top_{n} *_{n, n} x \in R_{2 n}$.
ii- $\mathbb{A}(R)={\underset{\longrightarrow}{\lim }}_{n \in \mathbb{N}} R_{n}$ is a boolean ring (or, equivalently, $\mathbb{T}(\mathbb{A}(R)) \in \operatorname{Igr}_{+}$).
Proof.
i- The property is clear if $n=0$. If $n \geq 1$, then the property can be verified by induction on the number of generators $k \geq 1, x=\sum_{i=1}^{k} a_{1, i} *_{1,1} a_{2, i} *_{1,1} \cdots *_{1,1} a_{n, i} \in R_{n}$ : if $k=1$, then note that

$$
\begin{aligned}
x *_{n, n} x & =\left(a_{1} * a_{2} * \cdots * a_{n}\right) *\left(a_{1} * a_{2} * \cdots * a_{n}\right) \\
& =\left(a_{1} * a_{1}\right) *\left(a_{2} * a_{2}\right) * \cdots\left(a_{n} * a_{n}\right)=\left(\top_{1} * a_{1}\right) *\left(\top_{1} * a_{2}\right) * \cdots *\left(\top_{1} * a_{n}\right) \\
& =\left(\top_{n}\right) *\left(a_{1} * a_{2} * \cdots * a_{n}\right) ;
\end{aligned}
$$

if $k>1$, write $x=y+z$, where $y, z \in R_{n}$ are have $<k$ generator and then, by induction,

$$
\begin{aligned}
x *_{n, n} x & =(y+z) *_{n, n}(y+z)=y *_{n, n} y+y *_{n, n} z+z *_{n, n} y+z *_{n, n} z \\
& =y *_{n, n} y+z *_{n, n} z=\mathrm{T}_{n} *_{n, n} y+\mathrm{T}_{n} *_{n, n} z \\
& =\mathrm{T}_{n} *_{n, n}(y+z)=\mathrm{T}_{n} *_{n, n} x
\end{aligned}
$$

ii- This follows directly from item (i) and the definition of the ring structure in $\mathbb{A}(R)=\underline{\lim }_{n \in \mathbb{N}} R_{n}$.

By the previous Proposition and the universal property of the boolean hull of an Igr (Definition 5.2.5(ii)), we obtain:
igr+co
Corollary 5.3.13. Let $R \in$ Igr $_{+}$. Then:
$i-\quad X_{\mathbb{T}(\mathbb{A}(R))} \approx X_{R}$.
ii- $\mathbb{A}(R) \cong B(R)$.

## Lemma 5.3.14.

$i$ - Given $R \in I g r_{\mathbb{1}}, S \in \operatorname{Igr}$ and $f: S \rightarrow j_{\mathbb{1}}(R)$, we have: $f$ is coordinatewise surjective iff $f_{1}: S_{1} \rightarrow R_{1}$ is a surjective morphism of pointed $\mathbb{F}_{2}$-modules.
ii - Given $R \in \operatorname{Igr} r_{\mathbb{1}}, S \in \operatorname{Igr}$ and $f, h \in \operatorname{Igr}\left(j_{\mathbb{1}}(R), S\right)$, we have $f=h$ if and only if $f_{1}=h_{1}$.
Let $R, S \in$ Igr. The inclusion function $\iota_{R}: \mathbb{1}(R) \rightarrow R$ and projection function $\pi_{R}: R \rightarrow \mathcal{Q}(R)$ induces respective natural transformations $\iota: \mathbb{1} \Rightarrow 1_{\text {Igr }}$ and $\pi: 1_{\text {Igr }} \Rightarrow \mathcal{Q}$. Moreover, we have a natural transformation can : $\mathcal{Q} \mathbb{1} \Rightarrow \mathbb{1} \mathcal{Q}$ given by the rule $\operatorname{can}_{n}\left(l\left(a_{1}\right) \ldots l\left(a_{n}\right)\right):=l\left(a_{1}\right) \ldots l\left(a_{n}\right)$, $n \geq 1$. ( $\operatorname{can}_{n}$ is well defined and is an isomorphism basically because both $\mathcal{Q} \mathbb{1}(R)$ and $\mathbb{1} \mathcal{Q}(R)$ are generated in level 1 by $R_{1}$ and both graded rings satisfies the relation $\left.l(a) l(-a)=0\right)$.

We have another immediate consequence of the previous results (and adjunctions):

## Lemma 5.3.15.

$i$ - For all $R \in \operatorname{Igr} r_{h}, \mathbb{1}(R) \in I g r_{+}$and can $R_{R}$ is an isomorphism.
ii - For all $R \in \operatorname{Igr} r_{\mathbb{1}}, \mathcal{Q}(R) \in I g r_{+}$and can $_{R}$ is an isomorphism.
iii - To get projective limits in Igr ${ }_{+}$is enough to restrict the projective limits obtained in Igr:

$$
\lim _{i \in I} R_{i} \cong \mathbb{1}\left(\lim _{\underset{i \in I}{ }} j_{+}\left(R_{i}\right)\right) .
$$

iv - To get inductive limits in Igr ${ }_{+}$is enough to restrict the inductive limits obtained in Igr:

$$
\lim _{i \in I} R_{i} \cong \mathcal{Q}\left(\lim _{i \in I} j_{+}\left(R_{i}\right)\right) .
$$

### 5.4 Examples and Constructions of Quadratic Interest

Definition 5.4.1. $A$ filtered ring is a tuple $A=\left(A,\left(J_{n}\right)_{n \geq 0},+, \cdot, 0,1\right)$ where:
$i-(A,+, \cdot, 0,1)$ is a commutative ring with unit.
ii - $J_{0}=A$ and for all $n \geq 1, J_{n} \subseteq A$ is an ideal.
iii - For all $n, m \geq 0, n \leq m \Rightarrow J_{n} \supseteq J_{m}$.
iv - For all $n, m \geq 0, J_{n} \cdot J_{m} \subseteq J_{n+m}$.
$v-J_{0} / J_{1} \cong \mathbb{F}_{2}$ (then $2=1+1 \in J_{1}$ ).
vi - For all $n \geq 0, J_{n} / J_{n+1}$ is a group of exponent 2 (then $2 \cdot J_{n} \subseteq J_{n+1}$ and $2^{n} \in J_{n}$ ).
A morphism $f: A \rightarrow A^{\prime}$ of filtered rings is a ring homomorphism such that $f\left(J_{n}\right) \subseteq J_{n}^{\prime}$. The category of filtered rings will be denoted by FRing.
gradfilt
Definition 5.4.2. We define the inductive graded ring associated functor

$$
\text { Grad : FRing } \rightarrow \text { Igr }
$$

for $f: \operatorname{FRing}(A, B)$ as follow: $\operatorname{Grad}(A):=\left(\left(\operatorname{Grad}(A)_{n}\right)_{n \geq 0},\left(t_{n}\right)_{n \geq 0}, *\right) \in \operatorname{Igr}$ is the igr where
$i$ - For all $n \geq 0, \operatorname{Grad}(A)_{n}:=\left(J_{n} / J_{n+1},+_{n}, 0_{n}, \top_{n}\right)$ is the exponent 2 group with distinguished element $\top_{n}:=2^{n}+J_{n+1}$.
ii- For all $n \geq 0, t_{n}: \operatorname{Grad}(A)_{n} \rightarrow \operatorname{Grad}(A)_{n+1}$ is defined by $t_{n}:=2 \cdot$-, i.e,

$$
\text { For all } a+J_{n+1} \in J_{n} / J_{n+1}, t_{n}\left(a+J_{n+1}\right):=2 \cdot a+J_{n+2} \in J_{n+1} / J_{n+2} \text {. }
$$

Observe that $t_{n}\left(\top_{n}\right)=\top_{n+1}$, i.e, $t_{n}$ is a morphism of pointed $\mathbb{F}_{2}$-modules.
iii- For all $n, m \geq 0$ the biadditive function $*_{n m}: \operatorname{Grad}(A)_{n} \times \operatorname{Grad}(A)_{m} \rightarrow \operatorname{Grad}(A)_{n+m}$ is defined by the rule

$$
\left(a_{n}+J_{n+1}\right) *_{m n}\left(b_{m}+J_{m+1}\right)=a_{n} \cdot b_{m}+J_{n+m+1} \in J_{n+m} / J_{n+m+1} .
$$

The group $A_{g}:=\bigoplus_{n \geq 0} \operatorname{Grad}(A)_{n}$ of exponent 2 and the induced application $*: A_{g} \times A_{g} \rightarrow A_{g}$ are such that $\left(A_{g}, *\right)$ is a commutative ring with unit $\mathrm{T}_{1}=\left(2+J_{2}\right) \in J_{1} / J_{2}$.
iv - For all $n \geq 1, t_{n}=\mathrm{T}_{1} *_{1 n-}$.

The morphism $\operatorname{Grad}(f) \in \operatorname{Igr}\left(\operatorname{Grad}(A), \operatorname{Grad}\left(A^{\prime}\right)\right)$ is defined by the following rules: for all $n \geq 0, f_{n}: \operatorname{Grad}(A)_{n} \rightarrow \operatorname{Grad}\left(A^{\prime}\right)_{n}$ is given by

$$
f_{n}\left(a+J_{n+1}\right):=f_{n}(a)+J_{n+1}^{\prime}
$$

Note that $f_{n}$ a homomorphism of $\mathbb{F}_{2}$-pointed modules and $\bigoplus_{n \geq 0} f_{n}:\left(A_{g}, *\right) \rightarrow\left(A_{g}^{\prime}, *\right)$ is a homomorphism of graded rings with unit.
igrcont
Definition 5.4.3. The functor of graded ring of continuous functions over a space $X$

$$
\mathcal{C}\left(X,_{-}\right): \operatorname{Igr} \rightarrow \operatorname{Igr}
$$

is the functor defined for $f: R \rightarrow S$ by

$$
i-\mathcal{C}(X, R)_{0}:=R_{0} \cong \mathbb{F}_{2}
$$

ii- for all $n \geq 1, \mathcal{C}(X, R)_{n}:=\mathcal{C}\left(X, R_{n}\right)$ as a pointed $\mathbb{F}_{2}$-module,
iii - for all $n, m \geq 0, *_{n m}^{X}: \mathcal{C}\left(X, R_{n}\right) \times \mathcal{C}\left(X, R_{m}\right) \rightarrow \mathcal{C}\left(X, R_{n+m}\right)$ is given by $\left(\alpha_{n}, \beta_{m}\right) \mapsto \alpha_{n} *_{n m}^{X} \beta_{m}$, where for $x \in X$,

$$
\alpha_{n} *_{n m}^{X} \beta_{m}(x)=\alpha_{n}(x) *_{n m} \beta_{m}(x) \in R_{n+m}
$$

iv- $\mathcal{C}(X, f)_{0}:=f_{0}$ as an homomorphism of pointed $\mathbb{F}_{2}$-modules $R_{0} \rightarrow S_{0}$.
$v-$ for all $n \geq 1, \mathcal{C}(X, f)_{n}:=\mathcal{C}\left(X, f_{n}\right):=f_{n} \circ_{-} \in p \mathbb{F}_{2}-\bmod \left(\mathcal{C}\left(X, R_{n}\right), \mathcal{C}\left(X, S_{n}\right)\right)$.
Remark 5.4.4. Let $X$ be a topological space and let $R \in I g r_{\mathbb{1}}$. Note that if $X$ is compact or $R \in I g r_{f i n}$, then $\mathcal{C}(X, R) \in I g r_{\mathbb{1}}$.
sgfilt
Definition 5.4.5. We define the continuous function filtered ring functor

$$
\mathcal{C}: S G \rightarrow \text { FRing }
$$

as follow: first, consider the functor $\mathcal{C}\left(X_{-}, \mathbb{Z}\right): S G \rightarrow$ Ring, composition of the (contravariant) functors "associated ordering space" $X_{-}: S G \rightarrow$ Topop and "continuous functions in $\mathbb{Z}$ ring" $\mathcal{C}(-, \mathbb{Z}):$ Top ${ }^{o p} \rightarrow$ Ring (here $\mathbb{Z}$ is endowed with the discrete topology).

Now we define the functor $\mathcal{C}: S G \rightarrow$ FRing: given a special group $G \in S G$, we define

$$
\mathcal{C}(G):=\left(R(G),\left(J_{n}(G)\right)_{n \geq 0},+, \cdot, 0,1\right)
$$

where
$i-(R(G),+, \cdot, 0,1)$ is the subring of $\mathcal{C}\left(X_{G}, \mathbb{Z}\right)$ of continuous functions of constant parity, i.e, $R(G):=J_{0}(G) \xrightarrow{i_{0}(G)} \mathcal{C}\left(X_{G}, \mathbb{Z}\right)$ is the image of the monomorphism of rings with unit

$$
j_{0}(G): \mathcal{C}\left(X_{G}, 2 \mathbb{Z}\right) \cup \mathcal{C}\left(X_{G}, 2 \mathbb{Z}+1\right) \rightarrow \mathcal{C}\left(X_{G}, \mathbb{Z}\right)
$$

ii - For all $n \geq 1, J_{n}(G) \xrightarrow{i_{n}(G)} J_{0}(G)$ is the ideal of $R(G)$ (and also of $\mathcal{C}\left(X_{G}, \mathbb{Z}\right)$ ) that is the image of the monomorphism of abelian groups

$$
j_{n}(G): \mathcal{C}\left(X_{G}, 2^{n} \mathbb{Z}\right) \rightarrow \mathcal{C}\left(X_{G}, 2 \mathbb{Z}\right) \cup \mathcal{C}\left(X_{G}, 2 \mathbb{Z}+1\right)
$$

We also have $J_{0}(G) / J_{1}(G) \cong \mathbb{F}_{2}$ and for all $n, m \geq 0$ :
$a$ - If $n \geq m$ then $J_{n}(G) \supseteq J_{m}(G)$;
$b-J_{n}(G) \cdot J_{m}(G) \subseteq J_{n+m}(G) ;$
$c-2 J_{n}(G)=J_{n+1}(G) \Rightarrow J_{n}(G) / J_{n+1}(G)$ is an exponent 2 group.
On the morphisms, for $f \in S G\left(G, G^{\prime}\right)$, we define $\mathcal{C}(f) \in \operatorname{FRing}\left(\mathcal{C}(G), \mathcal{C}\left(G^{\prime}\right)\right)$ by

$$
\mathcal{C}(f)(h)=\mathcal{C}\left(X_{f}, \mathbb{Z}\right)(h)
$$

for $h \in \mathcal{C}(G) . \mathcal{C}(f)$ is well-defined because $\mathcal{C}(f) \in \operatorname{Ring}\left(\mathcal{C}(G), \mathcal{C}\left(G^{\prime}\right)\right)$ and for all $n \geq 0$,

$$
\mathcal{C}(f)\left(J_{n}(G)\right) \subseteq J_{n}\left(G^{\prime}\right) .
$$

Definition 5.4.6. We define the continuous function graded ring functor by

$$
\text { Grad } \circ \mathcal{C}: S G \rightarrow \text { Igr } .
$$

For convenience, we describe this functor now: given $G \in \mathrm{SG}$,

$$
\operatorname{Grad}(\mathcal{C}(G)):=\left(\left(\operatorname{Grad}(\mathcal{C}(G))_{n}\right)_{n \geq 0},\left(t_{n}\right)_{n \geq 0}, \cdot\right)
$$

where:
i $-\operatorname{Grad}(\mathcal{C}(G))_{n}:=\left(J_{n}(G) / J_{n+1}(G), \cdot, 0 \cdot J_{n+1}(G), 2^{n} J_{n+1}(G)\right)$, where $2 \in \mathcal{C}\left(X_{G}, \mathbb{Z}\right)$ is the constant function of value $2 \in 2 \mathbb{Z} \subseteq \mathbb{Z}$.
ii - For all $n \geq 0, J_{n}(G) / J_{n+1}(G) \xrightarrow{t_{2}=2 \cdot} J_{n+1}(G) / J_{n+2}(G)$.
iii - For all $n, m \geq 0, *_{n m}: J_{n}(G) / J_{n+1}(G) \times J_{m}(G) / J_{m+1}(G) \rightarrow J_{n+m}(G) / J_{n+m+1}(G)$ is given by

$$
\left(h_{n}+J_{n+1}(G)\right) *_{n m}\left(k_{m}+J_{m+1}(G)\right)=h_{n} k_{m}+J_{n+m+1}(G) .
$$

On the morphisms, given $f \in S G\left(G, G^{\prime}\right)$, we have that

$$
\operatorname{Grad}(\mathcal{C}(f))=\left(\operatorname{Grad}(\mathcal{C}(f))_{n}\right)_{n \geq 0} \in \operatorname{Igr}\left(\operatorname { G r a d } \left(\mathcal{C}(G), \operatorname{Grad}\left(\mathcal{C}\left(G^{\prime}\right)\right),\right.\right.
$$

where for all $n \geq 0, \operatorname{Grad}(\mathcal{C}(f))_{n}: \operatorname{Grad}(\mathcal{C}(G))_{n} \rightarrow \operatorname{Grad}\left(\mathcal{C}\left(G^{\prime}\right)\right)_{n}$ is such that

$$
\operatorname{Grad}(\mathcal{C}(f))_{n}\left(h+J_{n+1}(G)\right)=\mathcal{C}(f)(h)+J_{n+1}^{\prime}\left(G^{\prime}\right) .
$$

## Proposition 5.4.7.

a-There is a natural isomorphism $\theta: G r a d \circ \mathcal{C} \xrightarrow{\cong} \mathbb{T} \circ \mathcal{C}\left(X_{-}, \mathbb{F}_{2}\right)$. In particular, for all $G \in S G$, $\operatorname{Grad}(\mathcal{C}(G)) \in \operatorname{Igr}_{+}$.
$b$ - For all $0<n \leq m<\omega, 2^{m-n} ._{-}: J_{n}(G) / J_{n+1}(G) \rightarrow J_{m} / J_{m+1}(G)$ is an isomorphism of groups of exponent 2 .
$c-$ For all $n \geq 1$, there is an isomorphism of groups of exponent 2

$$
\theta_{n}(G): J_{n}(G) / J_{n+1}(G) \stackrel{\cong}{\leftrightarrows} \mathcal{C}\left(X_{G}, \mathbb{F}_{2}\right),
$$

given by the rule

$$
\theta_{n}\left(h+J_{n}(G)\right)(\sigma):=h_{n}(\sigma) / 2^{n} \in \mathcal{C}\left(X_{G}, \mathbb{Z} / 2 \mathbb{Z}\right)
$$

$d$ - For all $0<n \leq m<\omega$ the following diagram commute:

filtered Witt ring functor
Definition 5.4.8. We define the filtered Witt ring functor

$$
\mathcal{W}: S G \rightarrow F R i n g
$$

for $f \in S G(G, H)$ as follow: given a special group $G \in S G$, we define

$$
\mathcal{W}(G):=\left(W(G), I^{n}(G)_{n \geq 0}, \oplus, \otimes,\langle \rangle,\langle 1\rangle\right)
$$

where for all $n \geq 0, I^{n}(G)$ is the $n$-th power of the fundamental ideal

$$
I(G):=\left\{\varphi \in W(G): \operatorname{dim}_{2}(\varphi)=0\right\}
$$

We define $\mathcal{W}(f) \in \operatorname{FRing}(\mathcal{W}(G), \mathcal{W}(H))$ by the rule $\mathcal{W}(f)(\varphi):=f \star \varphi$.
$\mathcal{W}(G)$ is a filtered commutative ring with unit because:
i- $(W(G), \oplus, \otimes,\langle \rangle,\langle 1\rangle) \in$ Ring.
ii - For all $n \geq 0, I^{n}(G) \subseteq W(G)$ is an ideal.
iii - For all $n, m \geq 0, n \leq m \Rightarrow I^{n}(G) \supseteq I^{m}(G)$.
iv - For all $n, m \geq 0, I^{n}(G) \otimes I(G) \subseteq I^{n+m}(G)$.
$\mathrm{v}-I^{0}(G):=W(G)$.
vi - $I^{0}(G) / I^{1}(G) \cong \mathbb{F}_{2}$.
vii - For all $n \geq 0,\left(I^{n}(G) / I^{n+1}(G), \oplus,\langle \rangle\right)$ is a group of exponent 2 with distinguished element $2^{n}+I^{n+1}(G)$, where $2^{n}=\otimes_{i<n}\langle 1,1\rangle$.

Definition 5.4.9. We define the graded Witt ring functor

$$
G r a d \circ \mathcal{W}: S G \rightarrow I g r .
$$

We register, again, the following result:
Proposition 5.4.10. For each $G \in S G$ we have $\operatorname{Grad}(\mathcal{W}(G)) \in I g r_{+}$.

For each commutative ring with unit $A$, we have

$$
t(A)=\{a \in A: \text { exists } n \geq 0 \text { with } n \cdot a=0\} \subseteq A
$$

is an ideal (the torsion ideal of $A$ ). The association $A \mapsto A / t(A)$ is the component on the objects of an endofunctor of Ring.

For each $G \in S G$ we have a ring homomorphism with unit $\operatorname{sgn}_{G}: W(G) \rightarrow \mathcal{C}\left(X_{G}, \mathbb{Z}\right)$ given by the rule

$$
\operatorname{sgn}_{G}\left(\left\langle a_{0}, \ldots, a_{n-1}\right\rangle\right)(\sigma):=\sum_{i=0}^{n-1} \sigma\left(a_{i}\right) .
$$

The Pfister Local-Global principle says that $\operatorname{sgn}_{G}$ induces a monomorphism

$$
\operatorname{rsgn}_{G}: W(G) / t(W(G)) \rightarrow \mathcal{C}\left(X_{G}, \mathbb{Z}\right) .
$$

For each $G \in S G$ we have $\operatorname{sgn}_{G}(W(G)) \subseteq \mathcal{C}\left(X_{G}, 2 \mathbb{Z}\right) \cup \mathcal{C}\left(X_{G}, 2 \mathbb{Z}+1\right)$ (since the signatures of classes of forms has the same parity of its dimension) and for all $n \geq 1, \operatorname{sgn}_{G}\left(I^{n}(G)\right) \subseteq \mathcal{C}\left(X_{G}, 2^{n} \mathbb{Z}\right)$ (since $I^{n}(G)$ is the abelian subgroup of $W(G)$ generated by classes of Pfister forms of dimension $2^{n}$ ).
$\operatorname{sgn}: \mathcal{W} \rightarrow \mathcal{C}$ (respectively $\operatorname{rsgn}: \mathcal{W} / t(\mathcal{W}) \rightarrow \mathcal{C}$ ) is the natural transformation between functors

$$
S G \underset{\mathcal{C}}{\stackrel{\mathcal{W}}{\Longrightarrow}} \text { FRing }
$$

that provide natural transformations between functors $S G \Longrightarrow$ Igr :

$$
\begin{aligned}
\text { Grad } \cdot \operatorname{sgn}: \operatorname{Grad} & \circ \mathcal{W} \rightarrow \operatorname{Grad} \circ \mathcal{C}, \text { respectively } \\
\operatorname{Grad} \cdot \operatorname{rsgn} & : \operatorname{Grad} \circ(\mathcal{W} / t(\mathcal{W})) \rightarrow \operatorname{Grad} \circ \mathcal{C} .
\end{aligned}
$$

Remember that [MC] ([LC]) and [WMC] ([WLC]) are conjectures about these natural transformations.
$\mathcal{C}$ is a particular case of $\mathcal{W}$ in the following sense: $\mathcal{C}: S G \rightarrow$ FRing is naturally isomorphic to the composition of functors $S G \xrightarrow{\gamma \circ \beta} S G \xrightarrow{\mathcal{W}}$ FRing.

### 5.5 The adjunction between PSG and $\operatorname{Igr}_{h}$

By the very definition of the K-theory of hyperfields (with the notations in Theorem 4.3.3) we define the following functor.

Definition 5.5.1 (K-theories Functors). With the notations of Theorem 4.3.3 we have a functors $k: \mathcal{H M F} \rightarrow I g r_{+}, k: \mathcal{P S M F} \rightarrow I g r_{+}$induced by the reduced $K$-theory for hyperfields.

Now, let $R \in \operatorname{Igr}_{h}$. We define a hyperfield $(\Gamma(R),+,-. \cdot, 0,1)$ by the following: firstly, fix an exponential isomorphism $e_{R}:\left(R_{1},+_{1}, 0_{1}, \top_{1}\right) \rightarrow(G(R), \cdot, 1,-1)$ (in agreement with Definition 5.3.10). This isomorphism makes, for example, an element $a *_{11}\left(T_{1}+b\right) \in R_{2}, a, b \in R_{1}$ take the form $\left(l_{R}(x)\right) *_{11}\left(l_{R}((-1) \cdot y)\right) \in R_{2}, x, y \in G(R)$. By an abuse of notation, we simply write $l_{R}(x) l_{R}(-y) \in R_{2}, x, y \in G(R)$. In this sense, an element in $Q_{2}$ has the form $l_{R}(x) l_{R}(-x)$, $x \in \Gamma(R)$, and we can extend this terminology for all $Q_{n}, n \geq 2$ (see Definition 5.3.4, and Lemma 5.3.11.

Now, let $\Gamma(R):=G(R) \cup\{0\}$ and for $a, b \in \Gamma(R)$ we define

$$
\begin{aligned}
&-a:=(-1) \cdot a, \\
& a \cdot 0=0 \cdot a:=0, \\
& a+0=0+a=\{a\}, \\
& a+(-a)=\Gamma(R), \\
& \text { for } a, b \neq 0, a \neq-b \text { define } \\
& a+b:=\{c \in \Gamma(R): \text { there exist } d \in G(R) \text { such that } \\
&\left.a \cdot b=c \cdot d \in G(R) \text { and } l_{R}(a) l_{R}(b)=l_{R}(c) l_{R}(d) \in R_{2}\right\} .
\end{aligned}
$$


(3per)
prespechf
Proposition 5.5.2. With the above rules, $(\Gamma(R),+,-.,, 0,1)$ is a pre-special hyperfield.
Proof. We will verify the conditions of Definition 1.2.7. Note that by the definition of multivalued sum once we proof that $\Gamma(R)$ is an hyperfield, it will be hyperbolic. In order to prove that $(\Gamma(R),+,-\cdot, 0,1)$ is a multigroup we follow the steps below. Here we use freely the properties in Lemma 5.3.11.
i - Commutativity and $(a \in b+0) \Leftrightarrow(a=b)$ are direct consequence of the definition of multivaluated sum and the fact that $l_{R}(a) l_{R}(b)=l_{R}(b) l_{R}(a)$.
ii - We will prove that if $c \in a+b$, then $a \in c-b$ and $b \in c-a$.
If $a=0$ (or $b=0$ ) or $a=-b$, then $c \in a+b$ means $c=a$ or $c \in a-a$. In both cases we get $a \in c-b$ and $b \in c-a$.

Now suppose $a, b \neq 0$ with $a \neq-b$. Let $c \in a+b$. Then $a \cdot b=c \cdot d$ and $l_{R}(a) l_{R}(b)=$ $l_{R}(c) l_{R}(d) \in R_{2}$ for some $d \in G(R)$. Since $G(R)$ is a multiplicative group of exponent 2, we have $a \cdot d=b \cdot c$ (and hence $a \cdot(-d)=c \cdot(-b))$. Note that

$$
\begin{aligned}
l_{R}(a) l_{R}(-d) & =l_{R}(a) l_{R}(-a b c)=l_{R}(a) l_{R}(b c)=l_{R}(a) l_{R}(b)+l_{R}(a) l_{R}(c) \\
& =l_{R}(c) l_{R}(d)+l_{R}(a) l_{R}(c)=l_{R}(c) l_{R}(d)+l_{R}(c) l_{R}(a)=l_{R}(c) l_{R}(a d) .
\end{aligned}
$$

Similarly,

$$
\begin{aligned}
l_{R}(b) l_{R}(-c) & =l_{R}(b) l_{R}(-a b d)=l_{R}(b) l_{R}(a d)=l_{R}(b) l_{R}(a)+l_{R}(b) l_{R}(d) \\
& =l_{R}(a) l_{R}(b)+l_{R}(b) l_{R}(d)=l_{R}(c) l_{R}(d)+l_{R}(b) l_{R}(d) \\
& =l_{R}(b c) l_{R}(d)=l_{R}(a d) l_{R}(d) .
\end{aligned}
$$

Then

$$
\begin{aligned}
l_{R}(a) l_{R}(-d)-l_{R}(b) l_{R}(-c) & =l_{R}(c) l_{R}(a d)-l_{R}(a d) l_{R}(d)= \\
& =l_{R}(c) l_{R}(a d)-l_{R}(d) l_{R}(a d)=l_{R}(-c d) l_{R}(a d) .
\end{aligned}
$$

But

$$
\begin{aligned}
l_{R}(-c d) l_{R}(a d) & =l_{R}(-c d) l_{R}(a)+l_{R}(-c d) l_{R}(d)= \\
& =l_{R}(-c d) l_{R}(a)+l_{R}(c) l_{R}(d)=l_{R}(a) l_{R}(-c d)+l_{R}(a) l_{R}(b) \\
& =l_{R}(a) l_{R}(-b c d)=l_{R}(a) l_{R}(-a)=0 .
\end{aligned}
$$

Then

$$
l_{R}(a) l_{R}(-d)=l_{R}(b) l_{R}(-c),
$$

proving that $a \in b-c$. Similarly we prove that $b \in-c+a$.
iii - Since $(G(R), \cdot, 1)$ is an abelian group, we conclude that $(\Gamma(R), \cdot, 1)$ is a commutative monoid. Beyond this, every nonzero element $a \in \Gamma(R)$ is such that $a^{2}=1$.
iv - $a \cdot 0=0$ for all $a \in \Gamma(R)$ is direct from definition.
v - For the distributive property, let $a, b, d \in \Gamma(R)$ and consider $x \in d(a+b)$. We need to prove that

$$
\begin{equation*}
x \in d \cdot a+d \cdot b . \tag{}
\end{equation*}
$$

It is the case if $0 \in\{a, b, d\}$ or if $b=-a$. Now suppose $a, b, d \neq 0$ with $b \neq-a$. Then there exist $y \in G(R)$ such that $x=d y$ and $y \in a+b$. Moreover, there exist some $z \in G(R)$ such that $y \cdot z=a \cdot b$ and $l_{R}(y) l_{R}(z)=l_{R}(a) l_{R}(b)$.
If $0 \in\{a, b, d\}$ or if $b=-a$ there is nothing to prove. Now suppose $a, b, d \neq 0$ with $b \neq-a$. Therefore $(d y) \cdot(d z)=(d a) \cdot(d b)$ and

$$
\begin{aligned}
l_{R}(d y) l_{R}(d z) & =l_{R}(d) l_{R}(d)+l_{R}(d) l_{R}(z)+l_{R}(d) l_{R}(y)+l_{R}(y) l_{R}(z) \\
& =l_{R}(d) l_{R}(d)+l_{R}(d)\left[l_{R}(z)+l_{R}(y)\right]+l_{R}(y) l_{R}(z) \\
& =l_{R}(d) l_{R}(d)+l_{R}(d) l_{R}(y z)+l_{R}(y) l_{R}(z) \\
& =l_{R}(d) l_{R}(d)+l_{R}(d) l_{R}(a b)+l_{R}(a) l_{R}(b) \\
& =l_{R}(d) l_{R}(d)+l_{R}(d) l_{R}(a)+l_{R}(d) l_{R}(b)+l_{R}(a) l_{R}(b) \\
& =l_{R}(d a) l_{R}(d b),
\end{aligned}
$$

so $l_{R}(d y) l_{R}(d z)=l_{R}(d a) l_{R}(d b)$. Hence we have $x=d y \in d \cdot a+d \cdot b$.
vi - Using distributivity we have that for all $a, b, c, d \in \Gamma(R)$

$$
d[(a+b)+c]=(d a+d b)+d c \text { and } d[a+(b+c)]=d a+(d b+d c)
$$

In fact, if $x \in(a+b)+c$, then $x \in y+c$ for $y \in a+b$. Hence

$$
d x \in d y+d c \subseteq d(a+b)+d c=(d a+d b)+d c
$$

Conversely, if $z \in(d a+d b)+d c$, then $z=w+d c$, for some $w \in d a+d b=d(a+b)$. But in this case, $w=d t$ for some $t \in a+b$. Then

$$
z \in d t+d c=d[t+c] \subseteq d[(a+b)+c]
$$

Similarly we prove that $d[a+(b+c)]=d a+(d b+d c)$.
vii - Let $a \in \Gamma(R)$ and $x, y \in 1-a$. If $a=0$ or $a=1$ then we automatically have $x \cdot y \in 1-a$, so let $a \neq 0$ and $a \neq 1$. Then $x, y \in G(R)$ and there exist $p, q \in \Gamma(R)$ such that

$$
\begin{aligned}
& x \cdot p=1 \cdot a \text { and } l_{R}(x) l_{R}(p)=l_{R}(1) l_{R}(a)=0 \\
& y \cdot q=1 \cdot a \text { and } l_{R}(y) l_{R}(q)=l_{R}(1) l_{R}(a)=0 .
\end{aligned}
$$

Then $(x y) \cdot(p q a)=1 \cdot a$ and

$$
\begin{aligned}
l_{R}(x y) l_{R}(p q a) & =l_{R}(x y) l_{R}(p)+l_{R}(x y) l_{R}(q)+l_{R}(x y) l_{R}(a) \\
& =l_{R}(y) l_{R}(p)+l_{R}(x) l_{R}(q)+l_{R}(x) l_{R}(a)+l_{R}(y) l_{R}(a) \\
& =l_{R}(y) l_{R}(p a)+l_{R}(x) l_{R}(q a) \\
& =l_{R}(y) l_{R}(x)+l_{R}(x) l_{R}(y)=0 .
\end{aligned}
$$

Then $x y \in 1-a$, proving that $(1-a)(1-a) \subseteq(1-a)$. In particular, since $1 \in 1-a$, we have $(1-a)(1-a)=(1-a)$.
viii - Finally, to prove associativity, we use Theorem 2.2 .8 Let $\langle a, b\rangle \equiv\langle c, d\rangle$ the relation defined for $a, b, c, d \in \Gamma(R) \backslash\{0\}$ by

$$
\langle a, b\rangle \equiv\langle c, d\rangle \text { iff } a b=c d \text { and } l_{R}(a) l_{R}(b)=l_{R}(c) l_{R}(d) .
$$

For $0 \notin\{a, b, c, d\}, a \neq-b$ and $a b=c d$, we have

$$
a+b=c+d \text { iff }\langle a, b\rangle \equiv\langle c, d\rangle .
$$

Using items (i)-(vii) we get that $(\Gamma(R) \backslash\{0\}, \equiv, 1,-1)$ is a pre-special group. Then by Theorem 2.2.8 we have that $M(\Gamma(R) \backslash\{0\}) \cong \Gamma(R)$ is a pre-special hyperfield, and in particular, ( $\Gamma(R)$ is associative.

Definition 5.5.3. With the notations of Proposition 5.5.2 we have a functor $\Gamma:$ Igr $_{+} \rightarrow$ PSMF defined by the following rules: for $R \in \operatorname{Igr} r_{+}, \Gamma(R)$ is the special hyperfield obtained in Proposition 5 5.5.2 and for $f \in \operatorname{Igr}_{+}(R, S), \Gamma(f): \Gamma(R) \rightarrow \Gamma(S)$ is the unique morphism such that the following diagram commute


In other words, for $x \in R$ we have

$$
\Gamma(f)(x)=\left(e_{S} \circ f_{1} \circ l_{R}\right)(x)=e_{S}\left(f_{1}\left(l_{R}(x)\right)\right) .
$$

Theorem 5.5.4. The functor $k: \mathcal{P S M F} \rightarrow I g r_{+}$is the left adjoint of $\Gamma:$ Igr $_{+} \rightarrow \mathcal{P S \mathcal { M F }}$. The unity of the adjoint is the natural transformation $\phi: 1_{\mathcal{P S M \mathcal { F }}} \rightarrow \Gamma \circ k$ defined for $F \in \mathcal{P S M F}$ by $\phi_{F}=e_{k(F)} \circ \rho_{F}$.

Proof. We show that for all $f \in \mathcal{P S M \mathcal { M }}(F, \Gamma(R))$ there is an unique $f^{\sharp}: \operatorname{Igr}_{+}(k(F), R)$ such that $\Gamma\left(f^{\sharp}\right) \circ \phi_{F}=f$. Note that $\phi_{F}=e_{k(F)} \circ \rho_{F}$ is a group isomorphism (because $e_{k(F)}$ and $\rho_{F}$ are group isomorphisms).

Let $f_{0}^{\sharp}: 1_{\mathbb{F}_{2}}: \mathbb{F}_{2} \rightarrow \mathbb{F}_{2}$ and $f_{1}^{\sharp}:=l_{R} \circ f \circ\left(\phi_{F}\right)^{-1} \circ e_{k(F)}: k_{1}(F) \rightarrow R_{1}$. For $n \geq 2$, define
$h_{n}: \prod_{i=1}^{n} k_{1}(F) \rightarrow R_{n}$ by the rule

$$
h_{n}\left(\rho\left(a_{1}\right), \ldots, \rho\left(a_{n}\right)\right):=l_{R}\left(f\left(a_{1}\right)\right) * \ldots * l_{R}\left(f\left(a_{n}\right)\right)
$$

We have that $h_{n}$ is multilinear and by the Universal Property of tensor products we have an induced morphism $\bigotimes_{i=1}^{n} k_{n}(F) \rightarrow R_{n}$ defined on the generators by

$$
h_{n}\left(\rho\left(a_{1}\right) \otimes \ldots \otimes \rho\left(a_{n}\right)\right):=l_{R}\left(f\left(a_{1}\right)\right) * \ldots * l_{R}\left(f\left(a_{n}\right)\right) .
$$

Now let $\eta \in Q_{n}(F)$. Suppose without loss of generalities that $\eta=\rho\left(a_{1}\right) \otimes \ldots \otimes \rho\left(a_{n}\right)$ with $a_{1} \in 1-a_{2}$. Then $f\left(a_{1}\right) \in 1-f\left(a_{2}\right)$ which imply $l_{R}\left(f\left(a_{1}\right)\right) \in 1-l_{R}\left(f\left(a_{2}\right)\right)$. Since $R_{n} \in \operatorname{Igr}_{+}$,

$$
h_{n}(\eta):=h_{n}\left(\rho\left(a_{1}\right) \otimes \ldots \otimes \rho\left(a_{n}\right)\right)=l_{R}\left(f\left(a_{1}\right)\right) * \ldots * l_{R}\left(f\left(a_{n}\right)\right)=0 \in R_{n}
$$

Then $h_{n}$ factors through $Q_{n}$, and we have an induced morphism $\bar{h}_{n}: k_{n}(F) \rightarrow R_{n}$. We set $f_{n}^{\sharp}:=\bar{h}_{n}$. In other words, $f_{n}^{\sharp}$ is defined on the generators by

$$
f_{n}^{\sharp}\left(\rho\left(a_{1}\right) \ldots \rho\left(a_{n}\right)\right):=l_{R}\left(f\left(a_{1}\right)\right) * \ldots * l_{R}\left(f\left(a_{n}\right) .\right.
$$

Finally, we have

$$
\begin{aligned}
\Gamma\left(f^{\sharp}\right) \circ \phi_{F} & =\left[e_{R} \circ\left(f_{1}^{\sharp}\right) \circ e_{k(F)}^{-1}\right] \circ\left[e_{k(F)} \circ \rho_{F}\right]=e_{R} \circ\left(f_{1}^{\sharp}\right) \circ \rho_{F} \\
& =e_{R} \circ\left[l_{R} \circ f \circ\left(\phi_{F}\right)^{-1} \circ e_{k(F)}\right] \circ \rho_{F} \\
& =f \circ\left(\phi_{F}\right)^{-1} \circ\left[e_{k(F)} \circ \rho_{F}\right] \\
& =f \circ\left(\phi_{F}\right)^{-1} \circ \phi_{F}=f .
\end{aligned}
$$

For the unicity, let $u, v \in \operatorname{Igr}_{+}(k(F), R)$ such that $\Gamma(u) \circ \phi_{F}=\Gamma(v) \circ \phi_{F}$. Since $\phi_{F}$ is an isomorphism we have $u_{1}=v_{1}$ and since $k(F) \in \operatorname{Igr}_{+}$we have $u=v$.

As we have already seen in Theorem 5.5.4, there natural transformation $\phi_{F}: F \rightarrow \Gamma(k(F))$ is a group isomorphism. Now let $a, c, d \in F$ with $a \in c+d$. Then $\phi_{F}(a) \in \phi_{F}(c)+\phi_{F}(d)$, i.e, $\phi_{F}$ is a morphism of hyperfields. In fact, if $0 \in\{a, c, d\}$ there is nothing to prove. Let $0 \notin\{a, c, d\}$. To prove that $\phi_{F}(a) \in \phi_{F}(c)+\phi_{F}(d)$ we need to show that $\rho_{F}(a) \rho_{F}(a c d)=\rho_{F}(c) \rho_{F}(d)$. In fact, from $a \in c+d$ we get $a c \in 1+a d$, and then $\rho_{F}(a c) \rho_{F}(a d)=0$. Moreover

$$
\begin{aligned}
\rho_{F}(a) \rho_{F}(a c d)+\rho_{F}(c) \rho_{F}(d) & =\rho_{F}(a) \rho_{F}(a c d)+\rho_{F}(c) \rho_{F}(d)+\rho_{F}(a c) \rho_{F}(a d) \\
& =\rho_{F}(a) \rho_{F}(a c)+\rho_{F}(a) \rho_{F}(d)+\rho_{F}(c) \rho_{F}(d)+\rho_{F}(a c) \rho_{F}(a d) \\
& =\left[\rho_{F}(a) \rho_{F}(a c)+\rho_{F}(a c) \rho_{F}(a d)\right]+\left[\rho_{F}(a) \rho_{F}(d)+\rho_{F}(c) \rho_{F}(d)\right] \\
& =\rho_{F}(d) \rho_{F}(a c)+\rho_{F}(d) \rho_{F}(a c)=0
\end{aligned}
$$

proving that $\phi_{F}(a) \in \phi_{F}(c)+\phi_{F}(d)$. Unfortunately we do not now if or where $\phi_{F}$ is a strong morphism. Then we propose the following definition.
kstable-def
Definition 5.5.5 (The $k$ stability). Let $F$ be a pre-special hyperfield. We say that $F$ is $k$-stable if $\phi_{F}: F \rightarrow \Gamma(F(G))$ is a full morphism. Alternatively, $F$ is $k$-stable if for all $a, b, c, d \in \dot{F}$, if $a b=c d$ then

$$
\rho_{F}(a) \rho_{f}(b)=\rho_{F}(c) \rho_{F}(d) \text { imply } a c \in 1+c d
$$

Proposition 5.5.6. Every PSG $G$ has a $k$-stable hull $G_{(k)}$ that satisfies the corresponding universal
property. This is just given by

$$
G_{(k)}=\lim _{n \in \mathbb{N}}(\Gamma \circ k)^{n}(G) .
$$

Thus the inclusion functor $P S G_{(k)} \hookrightarrow P S G$ has a left adjoint $(k): P S G \rightarrow P S G_{(k)}$.
We emphasize that if $G$ is $A P(3)$ special group, then $G$ is $k$-stable. In particular, every reduced special group is $k$-stable, and if $F$ is a field of characteristic not 2 , then $G(F)$ is also $k$-stable.

In the next Chapter, it is established the Arason-Pfister Hauptsatz (Theorem 6.3.2) for every special group $G$, (i.e., $G$ satisfies $A P(n)$ for each $n \in \mathbb{N}$.)

## Proposition 5.5.7.

$i$ - For each $G \in S G, \Gamma\left(s_{G}\right): \Gamma(\mathcal{K}(G)) \rightarrow \Gamma(\operatorname{Grad}(\mathcal{W}(G)))$ is a PSG-isomorphism.
ii - For each $G \in \mathcal{R S G}, \kappa_{G}: G \rightarrow \Gamma(\mathcal{K}(G))$ is a PSG-isomorphism.
iii - For each $G \in \mathcal{R S G}, \omega_{G}: G \rightarrow \Gamma(\operatorname{Grad}(\mathcal{W}(G)))$ is a PSG-isomorphism.
Proposition 5.5.8. Let $G$ be a PSG. Are equivalent:

$$
\begin{aligned}
& i-G \in \mathcal{P S} \mathcal{G}_{\text {fin }} . \\
& i i-\mathcal{K}(G) \in \operatorname{Igr}_{f i n} .
\end{aligned}
$$

Proposition 5.5.9. Let $G$ be a $S G$. Are equivalent:

$$
\begin{aligned}
& i-G \in S G_{f i n} \\
& i i-\mathcal{K}(G) \in I g r_{f i n} \\
& i i i-(\operatorname{Grad} \circ \mathcal{W})(G) \in \operatorname{Igr}_{f i n} .
\end{aligned}
$$

Proposition 5.5.10. The canonical arrow

$$
\text { can: } \underset{i \in I}{\lim } \mathcal{K}\left(G_{i}\right) \rightarrow \mathcal{K}\left(\underset{i \in I}{\lim } G_{i}\right)
$$

is an Igr $+_{+}$-isomorphism as long as the I-colimits above exists.
Proposition 5.5.11. The canonical arrow
is an Igr+-morphism pointwise surjective, as long as the I-colimits above exists.
Remark 5.5.12. In [27] there is an interesting analysis identifying the boolean hull of a special group $G$ (or special hyperfield $F=G \cup\{0\}$ ) with the boolean hull of the inductive graded rings $k_{*}(F), W_{*}(F) \in I g r_{+}$(see the above Corollary 5.3.13). It could be interesting to compare the space of orderings of $R \in$ Igr $_{h}$ and of $\Gamma(R) \in \mathcal{P S M \mathcal { F }}$.

### 5.6 Igr and Marshall's Conjecture

igrmarshall
Using the Boolean hull functor, M. Dickmann and F. Miraglia provide an encoding of Marshall's signature conjecture ([MC]) for reduced special groups by the condition

$$
\langle 1,1\rangle \otimes-: I^{n}(G) / I^{n+1}(G) \rightarrow I^{n+1}(G) / I^{n+2}(G)
$$

to be injective, for each $n \in \mathbb{N}$. In fact they introduce the notion of a [SMC] reduced special group:

$$
l(-1) \otimes-: k_{n}(G) \rightarrow k_{n+1}(G)
$$

is injective, for each $n \in \mathbb{N}$. They establish that, [SMC] imply [MC], for every reduced special group $G$. Moreover (see 5.1 and 5.4 in [30]):

- The inductive limit of [SMC] groups is [SMC].
- The finite product of [SMC] groups is [SMC].
- $G(F)$ is [SMC], for every Pythagorean field $F$ (with $(\operatorname{char}(F) \neq 2)$.


## Proposition 5.6.1.

$i-s: k \rightarrow G r a d \circ \mathcal{W}$ is a "surjective" natural transformation, where for each $G \in S G$ and all $n \geq 1, s_{n}(G): K_{n}(G) \rightarrow I^{n}(G) / I^{n+1}(G)$ is given by the rule

$$
s_{n}(G)\left(\sum_{i=0}^{s-1} l\left(g_{1, i}\right) \otimes \ldots \otimes l\left(g_{n, i}\right)+\mathcal{Q}_{n}(G)\right):=\bar{\bigotimes}_{i=0}^{s-1}\left[\left\langle 1,-g_{1, i}\right\rangle\right] \bar{\otimes} \ldots \bar{\otimes}\left[\left\langle 1,-g_{n, i}\right\rangle\right] \bar{\otimes} I^{n+1}(G)
$$

ii - $r: G r a d \circ \mathcal{W} \rightarrow k$ is a natural transformation, where for each $G \in S G$ and all $n \geq 1$, $r_{G}^{n}: I^{n}(G) / I^{n+2}(G) \rightarrow k_{2 n-1}(G)$ is given by the rule

$$
\begin{aligned}
r_{n}(G)( & \left.\bar{\bigotimes}_{i=0}^{s-1}\left[\left\langle 1,-g_{1, i}\right\rangle\right] \bar{\otimes} \ldots \bar{\otimes}\left[\left\langle 1,-g_{n, i}\right\rangle\right] \bar{\otimes} I^{n+1}(G)\right):= \\
& \sum_{i=0}^{s-1} l(-1)^{2^{n-1}-n} l\left(g_{1, i}\right) \otimes \ldots \otimes l\left(g_{n, i}\right)+\mathcal{Q}_{2 n-1}(G)
\end{aligned}
$$

iii - For all $n \geq 1, r_{n}(G) \circ s_{n}(G)=l(-1)^{2^{n-1}-n} \bar{\otimes}_{-}$.
iv- We have an isomorphism of pointed $\mathbb{F}_{2}$-modules: $s_{G}^{1}: k_{1}(G) \xrightarrow{\cong} I^{1}(G) / I^{2}(G), s_{G}^{2}: k_{2}(G) \xrightarrow{\cong}$ $I^{2}(G) / I^{3}(G)$.
$v$ - If $G$ is [SMC] Then $s_{G}: k(G) \rightarrow G r a d \circ \mathcal{W}(G)$ is an isomorphism.

We finish this chapter considering a general setting for "Marshall's conjectures", that includes the previous case of the Igr's $W_{*}(F), k_{*}(F)$ for special hyperfields $F$.

Let $R \in I g r_{+}$. The ideal, $\operatorname{nil}(R)$, in the $\operatorname{ring} \bigoplus_{n \in \mathbb{N}} R_{n}$, formed by all of its nilpotent elements, determines $N(R)$ a $\operatorname{Igr}$-ideal of $R$, where $(N(R))_{n}:=\operatorname{nil}(R) \cap R_{n}, \forall n \in \mathbb{N}$. Note that, by Proposition 5.3.12, $(\operatorname{nil}(R))_{n}=\left\{a \in R_{n}: \exists k \in \mathbb{N} \backslash\{0\}\left(\top_{k n} *_{k n, n} a=0_{(k+1) n}\right)\right\}, \forall n \in \mathbb{N}$.

Remark 5.6.2. Let $\rho: \mathbb{N} \rightarrow \mathbb{N}$ be an increasing function and define $\left(N_{\rho}(R)\right)_{n}=\left\{a \in R_{n}\right.$ : $\left.\exists k \in \mathbb{N}\left(\top_{\rho(n)} *_{\rho(n), n} a=0_{\rho(n)+n}\right)\right\}, \forall n \in \mathbb{N}$. Then $\left(N_{\rho}(R)\right)_{n}$ is a subgroup of $R_{n}$ and, since $\rho(n+k) \geq \rho(n)$, we have $\left(N_{\rho}(R)\right)_{n} *_{n, k} R_{k} \subseteq\left(N_{\rho}(R)\right)_{n+k}$. Summing up, $\left.\left(N_{\rho}(R)\right)_{n}\right)_{n \in \mathbb{N}}$ is an Igr-ideal.

The following result is straightforward consequence of the Definitions and 5.2.3, 5.3.13.
Proposition 5.6.3. For each $R \in I g r_{+}$are equivalent:
$i$ - For all $n \leq m \in \mathbb{N}, \operatorname{ker}\left(h_{n m}\right)=\left\{0_{n}\right\} \in R_{n}$.
ii - The canonical morphism $R \rightarrow \mathbb{T}(\mathbb{A}(R))$ is pointwise injective.
iii - There exists a boolean ring $B$ and a pointwise injective Igr-morphism $R \rightarrow \mathbb{T}(B)$.
Moreover, if $R \in I g r_{\text {fin }}$, these are equivalent to
$i v-N(R) \cong \mathbb{T}(0) \in \operatorname{Igr}$.
Motivated by item (i), we use the abbreviation $M C(R)$ to say that $R$ satisfies one (and hence all) of the above conditions.

In the following, we fix a category of $L$-structures $\mathcal{A}$ that is closed under directed inductive limits and a functor $F_{*}: \mathcal{A} \rightarrow I g r_{+}$be a functor that preserves directed inductive limits. Examples of such kind of functors are $k_{*}: \mathcal{H} \mathcal{M F} \rightarrow I g r_{+}$and $W_{*}: \mathcal{H} \mathcal{M F} \rightarrow I g r_{+}$, since such hyperfields can be conveniently described in the first-order relational language for multirings and it is closed under directed inductive limits. Related examples are the functors $k_{*}: S G \rightarrow I g r_{+}$and $W_{*}: S G \rightarrow I g r_{+}$; note that $S G$ is a full subcategory of $L_{S G}-S t r$ that is closed under directed inductive limits and under arbitrary products.

Proposition 5.6.4. If $(I, \leq)$ is an upward directed poset and $\Gamma:(I, \leq) \rightarrow \mathcal{A}$ is such that: $M C\left(F_{*}(\Gamma(i))\right)$, for all $i \in I$, then $M C\left(F_{*}\left(\lim _{i \in I} \Gamma(i)\right)\right)$.
Proof. The hypothesis on $F_{*}$ and the fact that the directed inductive limits in $I g r_{+}$are pointwise, give us immediately that the mappings $h_{n}: F_{n}\left(\lim _{i \in I} \Gamma(i)\right) \rightarrow F_{n+1}\left(\lim _{i \in I} \Gamma(i)\right)$ are isomorphic to the injective maps $\lim _{i \in I} h_{n}^{i}:{\underset{\longrightarrow}{\rightarrow}}_{\lim } F_{n}(\Gamma(i)) \rightarrow{\underset{\longrightarrow}{i}}_{i \in I} F_{n+1}(\Gamma(i))$, for each $n \in \mathbb{N}$. Therefore it holds

$$
M C\left(F_{*}(\underset{i \in I}{\lim } \Gamma(i))\right)
$$

Corollary 5.6.5. Let $F \subseteq P(I)$ be a filter and let $\left\{M_{i}: i \in I\right\}$ be a family of (non-empty) $L$ structures in $\mathcal{A}$. Suppose that $\mathcal{A}$ is closed under products and suppose that holds $\operatorname{MC}\left(F_{*}\left(\prod_{i \in J} M_{i}\right)\right)$, for each $J \in F$. Then holds $M C\left(F_{*}\left(\prod_{i \in J} M_{i} / F\right)\right)$.

Proof. This follows from the preceding result since, by a well-known model-theoretic result due to D. Ellerman ([34]), any reduced product of a family of (non-empty) $L$-structures, $\left\{M_{i}: i \in\right.$ $I\}$, module a filter $F \subseteq P(I)$, is canonically isomorphic to an upward directed inductive limit, ${\underset{\longrightarrow}{\longrightarrow}}_{J \in F}\left(\prod_{i \in J} M_{i}\right) \cong\left(\prod_{i \in I} M_{i}\right) / F$.

Proposition 5.6.6. Let $F_{*}: \mathcal{A} \rightarrow I g r_{+}$preserves pure embeddings. More precisely, if $M, M^{\prime} \in \mathcal{A}$ and $j: M \rightarrow M^{\prime}$ is a pure L-embedding, then $F_{*}(j): F_{*}(M) \rightarrow F_{*}\left(M^{\prime}\right)$ is a pure morphism of Igr's (described in the first-order polysorted language for Igr's).

Proof. This follows from the well known characterization result:
Fact: Let $L^{\prime}$ be a first-order language and $f: A \rightarrow B$ be an $L^{\prime}$-homomorphism. Then are equivalent

- $f: A \rightarrow B$ is a pure $L^{\prime}$-embedding.
- There exists an elementary $L^{\prime}$-embedding $e: A \rightarrow C$ and a $L^{\prime}$-homomorphism $h: B \rightarrow C$, such that $e=h \circ f$.
- There exists an ultrapower $A^{I} / U$ and a $L^{\prime}$-homomorphism $g: B \rightarrow A^{I} / U$, such that $\delta_{A}^{(I, U)}=$ $g \circ f$, where $\delta_{A}^{(I, U)}: A \rightarrow A^{I} / U$ is the diagonal (elementary) $L^{\prime}$-embedding.

Since the morphism $j: M \rightarrow M^{\prime}$ is a pure embedding, by the Fact there exists an ultrapower $M^{I} / U$ and a $L$-homomorphism $g: M^{\prime} \rightarrow M^{I} / U$, such that $\delta_{(I, U)}^{M}=g \circ j$, where $\delta_{M}^{(I, U)}: M \rightarrow M^{I} / U$ is the diagonal (elementary) $L$-embedding.

Since we have a canonical isomorphism can : $\lim _{\longrightarrow} \in \in U$ $M^{J} \xrightarrow{\cong} M^{I} / U$, applying the functor $F_{*}$, we obtain $F_{*}\left(M^{I} / U\right) \cong F_{*}\left(\lim _{\longrightarrow} J \in U\right.$ ( $\left.M^{J}\right) \cong \underline{\lim }_{\vec{J} \in U} \overrightarrow{F^{*}}\left(M^{J}\right) \rightarrow \underline{\lim }_{\vec{\longrightarrow}} J \in U\left(F^{*}(M)\right)^{J} \cong\left(F_{*}(M)\right)^{I} / U$.

Keeping track, we obtain that the above morphism $t: F_{*}\left(M^{I} / U\right) \rightarrow\left(F_{*}(M)\right)^{I} / U$ establishes a comparison between $F_{*}\left(\delta_{(I, U)}^{M}\right): F_{*}(M) \rightarrow F_{*}\left(M^{I} / U\right)$ and $\left.\delta_{(I, U)}^{F_{*}(M)}\right): F_{*}(M) \rightarrow F_{*}(M)^{I} / U$

$$
\left.\delta_{(I, U)}^{F_{*}(M)}\right)=t \circ F_{*}\left(\delta_{(I, U)}^{M}\right) .
$$

Since $F_{*}\left(\delta_{(I, U)}^{M}\right)=F_{*}(g) \circ F_{*}(j)$, combining the equations we obtain

$$
\left.\delta_{(I, U)}^{F_{*}(M)}\right)=t \circ F_{*}(g) \circ F_{*}(j) .
$$

Applying again the Fact, we conclude that $F_{*}(j): F_{*}(M) \rightarrow F_{*}\left(M^{\prime}\right)$ is a pure morphism of Igr's.

Corollary 5.6.7. For each $n \in \mathbb{N}$, the functor $F_{n}: \mathcal{A} \rightarrow p \mathbb{F}_{2}$ - mod preserves pure embeddings. More precisely, if $M, M^{\prime} \in \mathcal{A}$ and $j: M \rightarrow M^{\prime}$ is a pure L-embedding, then $F_{n}(j): F_{n}(M) \rightarrow$ $k_{n}\left(M^{\prime}\right)$ is a pure morphism of pointed $\mathbb{F}_{2}$-modules (described in the first-order single sorted language adequate). In particular $F_{n}(j): F_{n}(M) \rightarrow F_{n}\left(M^{\prime}\right)$ is an injective morphism of pointed $\mathbb{F}_{2}$-modules.

Corollary 5.6.8. Let $M, M^{\prime} \in \mathcal{A}$ and $j: M \rightarrow M^{\prime}$ is a pure L-embedding. If $M C\left(F_{*}\left(M^{\prime}\right)\right)$, then $M C\left(F_{*}(M)\right)$.

Proof. This follows directly from the previous Corollary. Indeed, suppose that holds $M C\left(F_{*}\left(M^{\prime}\right)\right)$. Since $h_{n}^{\prime}: F_{n}\left(M^{\prime}\right) \rightarrow F_{n+1}\left(M^{\prime}\right)$ and $F_{n}(j): F_{n}(M) \rightarrow F_{n}\left(M^{\prime}\right)$ are injective morphisms, then, by a diagram chase, $h_{n}: F_{n}(M) \rightarrow F_{n+1}(M)$ is an injective morphism too, thus holds $M C\left(F_{*}(M)\right)$.


## Chapter 6

## Quadratic Extensions of Special Groups, Hauptsatz and Consequences

In this Chapter we develop the theory of quadratic extensions for hyperfields/superfields, through the development of results concerning the superrings of polynomials, envisaging some applications to algebraic theory of quadratic forms and Real Algebraic Geometry. The main results here are the Arason-Pfister Hauptsatz for all special groups (Theorem 6.3.2) and its consequences.

The use of hyperfields/hyperrings/multirings in connection with Real Algebraic Geometry started 15 years ago, in 47].

The significance of these multivalued methods - as addition of roots to a superfield (Theorem 3.5.4) and Marshall's quotient of a superring (Theorem 6.1.5)-to (univalent) Commutative Algebra is indicated by applying these results to algebraic theory of quadratic forms: (i) obtaining new relevant constructions in the category of special groups (or its equivalent category special hyperfields, as in Theorems $6.2 .6,6.2 .12$; (ii) extending to all special hyperfields the validity of the Arason-Pfister Hauptsatz (Theorem 6.3.2)- a positive answer ([7]) to a question posed by Milnor in a classical paper of 1970 ([52], [7])- and established by Dickmann-Miraglia to the realm of reduced special groups (or its equivalent category real reduced hyperfields) in 2000 ([28]); and applied that to obtain interesting properties of graded rings associated to special hyperfields ([30, [18]).

Throughout this Chapter, all superrings will be considered associative.

### 6.1 Marshall's Quotient of Superfields

quotient-section
In the realm of multirings, the notion of the so called "Marshall's quotient", introduced in 47] and further developed in [24], is a quotient multiring defined for pair $(A, S)$ where $A$ is a multiring and $S \subseteq A$ is a multiplicative subset: given $a, b \in A$,

$$
a \approx_{S} b \text { iff there are } x, y \in S \text { such that } a x=b y
$$

Now we introduce the following:
Definition 6.1.1. Let $A$ be a superring and $S \subseteq A$. The set $S$ is called Marshall's coherent if it is multiplicative $(1 \in S$ and $S \cdot S \subseteq S$ ) and given $x, a \in A$ with $x \in$ as for some $s \in S$, there are $P, Q \subseteq S$ such that $x P=a Q$. We say that $S$ is nontrivial Marshall's coherent if $0 \notin S$.

Let $A$ be a superring with $S \subseteq A$ Marshall's coherent. For $a, b \in A$, define

$$
a \sim_{S} b \text { iff there are non-empty subsets } X, Y \subseteq S \text { with } a X=b Y
$$

Fact 6.1.2. If $A$ is a multiring viewed as a superring, then every multiplicative subset $S \subseteq A$ is Marshall's coherent and the above quotient notion coincides with the original Marshall's quotient, i.e. $\approx_{S}=\sim_{S}$.

## Lemma 6.1.3.

$i$ - For $a, b \in A$, the following are equivalent:
a) $a \sim b$.
b) There exists $s, t \in S$ such that as $\cap b t \neq \emptyset$.
c) There are $s, t, p, q \in S$ with $a(s t)=b(p q)$.
$i i$ - The relation $\sim$ is an equivalence relation.
Proof. we only need to deal with the case $S$ nontrivial.
i - The implication $c) \Rightarrow a$ ) is straightforward. For $a) \Rightarrow b$ ), let $X, Y \subseteq S$ such that $a X=b Y$. Then there are $s \in X$ and $t \in Y$ such that $a s \cap b t \neq \emptyset$. On the other hand, for $b) \Rightarrow c$ ), let $x \in a s \cap b t$. Thus, by Marshall's coherence, there are $M, N, P, Q \subseteq S$ such that $x M=a P$ and $x N=b Q$. Therefore,

$$
a(P N)=x(M N)=b(Q M) .
$$

ii - Let $a, b, c \in A$.

- Since $a \cdot\{1\}=a \cdot\{1\}$ and $1 \in S$, we have $a \sim a$.
- If $a \sim b$, then $a X=b Y$ for some $X, Y \subseteq S$. So $b Y=a X$ and $b \sim a$.
- Let $a \sim b$ and $b \sim c$. Then $a X=b Y$ and $b Z=c W$ for some $X, Y, Z, W \subseteq S$. Hence

$$
a(X Z)=b(Y Z)=c(W Y)
$$

and so $a \sim c$.

Now, let $A / m S$ be the set of equivalence classes of $\sim$. We want to prescribe a superring structure for $A /{ }_{m} S$.

For $a \in A$, let $[a]$ be the equivalence class of $a$ in $A /{ }_{m} S$. Define for $[a],[b] \in A /{ }_{m} S$ the congruence relations:

$$
\begin{aligned}
& {[c] \in[a]+[b] \text { iff there exist } c^{\prime}, a^{\prime}, b^{\prime} \in A \text { with } c^{\prime} \in a^{\prime}+b^{\prime} \text { and } c^{\prime} \sim c, a^{\prime} \sim a, b^{\prime} \sim b .} \\
& {[c] \in[a][b] \text { iff there exist } c^{\prime}, a^{\prime}, b^{\prime} \in A \text { with } c^{\prime} \in a^{\prime} \cdot b^{\prime} \text { and } c^{\prime} \sim c, a^{\prime} \sim a, b^{\prime} \sim b .} \\
& {[-a]:=-[a] .}
\end{aligned}
$$

lemsum1
Lemma 6.1.4. Let $A$ be a superring and $S \subseteq A$ a Marshall's coherent subset. Let $a, b, c \in A$.
$i-[c] \in[a]+[b]$ iff there is $s \in S$ such that $c s \subseteq a S+b S$.
ii - $[c] \in[a] \cdot[b]$ iff there is $s \in S$ such that $c s \subseteq a b S$.
Proof. we only need to deal with the case $S$ nontrivial.
i - $(\Rightarrow)$ : Let $c^{\prime}, a^{\prime}, b^{\prime} \in A$ such that $c^{\prime} \in a^{\prime}+b^{\prime}$ and $c^{\prime} \sim c, a^{\prime} \sim a, b^{\prime} \sim b$. Then

$$
c^{\prime} X^{\prime}=c X, a^{\prime} Y^{\prime}=a Y, b^{\prime} Z^{\prime}=b Z \text { for some } X, Y, Z, X^{\prime}, Y^{\prime}, Z^{\prime} \subseteq S
$$

and so

$$
c\left(X Y^{\prime} Z^{\prime}\right)=c^{\prime}\left(X^{\prime} Y^{\prime} Z^{\prime}\right) \in a^{\prime}\left(X^{\prime} Y^{\prime} Z^{\prime}\right)+b^{\prime}\left(X^{\prime} Y^{\prime} Z^{\prime}\right)=a\left(X^{\prime} Y Z^{\prime}\right)+b\left(X^{\prime} Y^{\prime} Z\right) \subseteq a S+b S
$$

Therefore, for any $s \in X Y^{\prime} Z^{\prime} \subseteq S$, we have $c s \subseteq a S+b S$.
$(\Leftarrow)$ : By hypothesis, there is $c^{\prime} \in c s \cap a t+b v$ for some $t, v \in S$. Therefore there exists $a^{\prime} \in a t$ and $b^{\prime} \in b v$ with $c^{\prime} \in a^{\prime}+b^{\prime}$. Lastly, Marshall's coherence implies that $c^{\prime} \sim c, a^{\prime} \sim a$ and $b^{\prime} \sim b$.
ii - $(\Rightarrow)$ : Let $c^{\prime}, a^{\prime}, b^{\prime} \in A$ such that $c^{\prime} \in a^{\prime} \cdot b^{\prime}$ and $c^{\prime} \sim c, a^{\prime} \sim a, b^{\prime} \sim b$. Then

$$
c^{\prime} X^{\prime}=c X, a^{\prime} Y^{\prime}=a Y, b^{\prime} Z^{\prime}=b Z \text { for some } X, Y, Z, X^{\prime}, Y^{\prime}, Z^{\prime} \subseteq S
$$

and so

$$
c\left(X Y^{\prime} Z^{\prime}\right)=c^{\prime}\left(X^{\prime} Y^{\prime} Z^{\prime}\right) \in a^{\prime}\left(X^{\prime} Y^{\prime} Z^{\prime}\right) b^{\prime}\left(X^{\prime} Y^{\prime} Z^{\prime}\right)=a\left(X^{\prime} Y Z^{\prime}\right) b\left(X^{\prime} Y^{\prime} Z\right) \subseteq b c S .
$$

Therefore, for any $s \in X Y^{\prime} Z^{\prime} \subseteq S$, we have $c s \subseteq a b S$.

Theorem 6.1.5. Let $A$ be a superring and $S \subseteq A$ a Marshall's coherent subset.
$i$ - The structure $(A / m S,+, \cdot,-,[0],[1])$ is a superring.
ii - The projection map $\pi: A \rightarrow A /{ }_{m} S$ is a universal morphism satisfying $\pi(S)=\{1\}$, that is, given a morphism $f: A \rightarrow B$ with $f(S)=\{1\}$, there is an unique morphism $\bar{f}: A /{ }_{m} S \rightarrow B$ such that $f=\bar{f} \circ \pi$. In other words, for every morphism $f: A \rightarrow B$ such that $f[S]=\{1\}$, there exist a unique morphism $\tilde{f}: A /{ }_{m} S \rightarrow B$ such that the following diagram commute:

where $\pi: A \rightarrow A /{ }_{m} S$ is the canonical projection $\pi(a)=\bar{a}$.
Proof. we only need to deal with the case $S$ nontrivial.
i - Firstly, we prove that $A /{ }_{m} S$ is a superring.

- $(A / m S,+,-,[0])$ is a multigroup.

The commutativity of + is straightforward. Let $a, b, c \in A$.
$-[c] \in[a]+[b] \Rightarrow-[a] \in-[b]+[c]$.
Let $c^{\prime}, a^{\prime}, b^{\prime} \in A$ with $c^{\prime} \in a^{\prime}+b^{\prime}$ and $c^{\prime} \sim c, a^{\prime} \sim a, b^{\prime} \sim b$. Then $-a^{\prime} \in-c^{\prime}+b^{\prime}$ and so $-[a] \in-[b]+[c]$.
$-[a]+[0]=\{[a]\}$.
Let $[x] \in[a]+[0]$. Then there is $s \in S$ such that $x s \subseteq a S$. So, by Marshall's coherence, $[x]=[a]$.
$-([a]+[b])+[c] \subseteq[a]+([b]+[c])$.
Let $[x] \in[y]+[c]$, with $[y] \in[a]+[b]$. Then there are $s, t \in S$ such that $x s \subseteq$ $y S+c S, y t \subseteq a S+b S$. Thus

$$
x s t \subseteq(a S+b S)+c S \subseteq a S+(b S+c S)
$$

It follows by Marshall's coherence that $[x] \in[a]+[l],[l] \in[b]+[c]$.

- $(A / m S, \cdot,[1])$ is a multimonoid.

The commutativity of $\cdot$ is straightforward too. So let $a, b, c \in A$.
$-([a] \cdot[b]) \cdot[c] \subseteq[a] \cdot([b] \cdot[c])$.
Let $[x] \in[y] \cdot[c]$ with $[y] \in[a][b]$. Then there are $s, t \in S$ such that $x s \subseteq y c S, y t \subseteq a b S$. Thus

$$
x s t \subseteq(a b) c S \subseteq a(b c) S
$$

and by Marshall's coherence $[x] \in[a][l],[l] \in[b][c]$.
$-[a] \cdot[1]=\{[a]\}$.
Let $[x] \in[a] \cdot[1]$. Then there is $s \in S$ with $x s \subseteq a S$. By Marshall's coherence $[x]=[a]$.

The verification of axioms $i i i, i v$ and $v$ of Definition 3.1.1 are straightforward.
ii - It follows immediately from Marshall's quotient definition that $\pi: A \rightarrow A /{ }_{m} S$ is a morphism satisfying $\pi(S)=\{1\}$. Now let $f: A \rightarrow B$ be a morphism with $f(S)=\{1\}$. Note that if $a \sim b$, then $a s=b t$ for some $s, t \in S$ and so for any $x \in a s \cap b t$ we have $f(x) \in f(a s) \cap f(b t) \subseteq$ $\{f(a)\} \cap\{f(b)\}$. Thus $f(a)=f(b)$. Then we can define $\bar{f}: A / m S \rightarrow B$ by $\bar{f}([a])=f(a)$. Since the multioperations are defined by congruence relations, we have that $f$ is a superring morphism.

As an immediate consequence, note that if $S, T$ are Marshall coherent subsets of $A$ and $S \subseteq T$, then we have a canonical surjective morphism of superrings:

$$
A /{ }_{m} S \rightarrow A /{ }_{m} T
$$

Corollary 6.1.6. Let $A$ be a superring and $S \subseteq A$ be a non-trivial Marshall coherent subset of $A$.
$i$ - If $A$ is full then $A / m S$ is full.
ii - If $A$ is a superdomain then $A /{ }_{m} S$ is a superdomain.
iii - If $A$ is a superfield then $A /{ }_{m} S$ is a superfield.

Theorem 6.1.7. Let $A$ be a superdomain and $S \subseteq A \backslash\{0\}$ such that $1 \in S, 0 \notin S, S \cdot S \subseteq S$ and $A^{2} \backslash\{0\} \subseteq S$. Then $S$ is a Marshall coherent subset of $S$. Moreover $A / m S$ is a hyperfield, i.e, for all $[a],[b] \in A /{ }_{m} S,[a][b]$ is a singleton set.

Proof. Let $a \in A, s \in S$ and $x \in a s$. Suppose without loss of generality that $x \neq 0$ (because if $x=0$ then $a=0$ because $A$ is a superdomain). Then $x a \subseteq a^{2} s \subseteq S$ and since

$$
x(x a)=a x^{2} \text { with } x a \text { and } x^{2} \text { contained in } S,
$$

we have that $S$ is a Marshall coherent subset. Now let $[c],[d] \in[a] .[b] \neq \emptyset$. Then $c s_{1} \subseteq a b S$ and $d s_{2} \subseteq a b S$ for some $s_{1}, s_{2} \in S$ (see Lemma 6.1.4).

If $c=0$ or $d=0$ then $0 \in a b S$ which imply $a=0$ or $b=0$ and $[a][b]=\{[0]\}$. Let $c, d \neq 0$. Then

$$
c d s_{1} s_{2} \subseteq a^{2} b^{2} S \cdot S \subseteq S
$$

Using this fact we get

$$
c\left(c d s_{1} s_{2}\right)=d\left(c^{2} s_{1} s_{2}\right) \text { with } c d s_{1} s_{2} \text { and } c^{2} s_{1} s_{2} \text { contained in } S .
$$

Moreover $c \sim d$, thus $[a] .[b]$ is a singleton.
We already know that $A /{ }_{m} S$ is a superdomain. To show that $A /{ }_{m} S$ is a superfield, it suffices to note that for each $a \in A$ such that $[a] \neq[0],[1] \in[a] .[a]$, or, equivalently, there is $s \in S$ such that $1 s \subseteq a S . a S$, but $a^{2} \in S$ and $1 . a^{2} \subseteq(a .1)$.(a.1). Thus $A /{ }_{m} S$ is superfield with single-valued products, i.e., $A /{ }_{m} S$ is an hyperfield.

From Theorem 6.1.7, we have many examples of Marshall coherent sets for superdomains $A$ (of course, after removing zero):

- the squares

$$
A^{2}:=\bigcup_{a \in A} a^{2} ;
$$

- the sum of squares

$$
\sum A^{2}:=\bigcup_{a_{1}, \ldots, a_{n} \in A, n \in \mathbb{N}} a_{1}^{2}+a_{2}^{2}+\ldots+a_{n}^{2}
$$

- preorderings, that are subsets $T \subseteq A$ with $T+T \subseteq T, T \cdot T \subseteq T$ and $A^{2} \subseteq T$.


### 6.2 Special Hyperfields and Quadratic Extensions

quadrapre-keasmien
Theorem 6.2.1. Let $F$ be a hyperbolic hyperfield such that $1+1=\{0,1\}$ and $1=-1$. Then $F \cong\{0,1\}$ (the Krasner hyperfield). In particular, $F$ is a DM-hyperfield.

Proof. Just observe that

$$
F=1-1=1+1=\{0,1\} .
$$

Throughout this section we establish the following notation: Let $G$ be a formally real pre-special group and $F$ its special multifield associated. In particular, if $\alpha \in F, \alpha \neq 0,1$, the polynomial $f(X) \in F[X], f(X)=X^{2}-\alpha$ has no roots in $F$ (basically because $\alpha^{2}=1$ for all $\alpha \in F^{*}$ ).

Then, let $\omega \in \bar{F}$, such that $0 \in f(\omega)$ (i.e, $\left.0 \in \omega^{2}-\alpha\right)$ and $\operatorname{Irr}(F, \omega)=f$. Note that this imply
$\alpha=\omega^{2}$. Next, consider the superfield extension $F(\omega)=F(\operatorname{Irr}(F, \omega))$ and let

$$
\begin{aligned}
S_{F}(\omega) & =(F(\omega)) / m\left(F(\omega)^{2} \backslash\{0\}\right) \\
S_{F}^{r e d}(\omega) & =(F(\omega)) / m\left(\sum F(\omega)^{2} \backslash\{0\}\right)
\end{aligned}
$$

We denote $\omega=\sqrt{\alpha}$.
prop1
Proposition 6.2.2. Let $F$ be a formally real special hyperfield and $S_{F}(\omega)$ as above.
$a$ - Let $a_{1}, a_{2}, \ldots, a_{n}, b_{1}, b_{2}, \ldots, b_{n} \in F$. Then

$$
\left(a_{1}+b_{1} \omega\right)^{2}+\left(a_{2}+b_{2} \omega\right)^{2}+\ldots+\left(a_{n}+b_{n} \omega\right)^{2} \subseteq(1+\alpha)+2\left[a_{1} b_{1}+a_{2} b_{2}+\ldots+a_{n} b_{n}\right] \omega,
$$

where $2 X:=X+X$.
$b-\sum F(\omega)^{2} \backslash\{0\}=(1+\alpha)+F \cdot \omega$.
$c$ - Denote (2) $F=\bigcup\{x+x: x \in F\}$. Then

$$
F(\omega)^{2}=(1+\alpha)+(2) F \cdot \omega .
$$

$d-F(\omega)^{2}=\sum F(\omega)^{2}$ iff (2)F=F.
$e--1 \notin \sum F(\omega)^{2}$ iff $-1 \notin 1+\alpha$.
$f-\omega \notin F(\omega)^{2}$.
$g$ - The morphism $F \rightarrow S_{F}(\omega)$ is full and not injective.
$h$ - If $-1 \notin 1+\alpha$ and $a, b \in a+b$ for all $a, b \in \dot{F}$ then $S_{F}(\omega)$ is a real reduced hyperfield (and then, a reduced special group).

## Proof.

a - In fact, if $n=2$, then

$$
\begin{aligned}
\left(a_{1}+b_{1} \omega\right)^{2}+\left(a_{2}+b_{2} \omega\right)^{2} & =a_{1}^{2}+a_{1} b_{1} \omega+a_{1} b_{1} \omega+b_{1}^{2} \omega^{2}+a_{2}^{2}+a_{2} b_{2} \omega+a_{2} b_{2} \omega+b_{2}^{2} \omega^{2} \\
& =1+\left[a_{1} b_{1}+a_{1} b_{1}\right] \omega+\alpha+1+\left[a_{2} b_{2}+a_{2} b_{2}\right] \omega+\alpha \\
& =[1+\alpha+1+\alpha]+\left[a_{1} b_{1}+a_{1} b_{1}+a_{2} b_{2}+a_{2} b_{2}\right] \omega \\
& \subseteq(1+\alpha)+2\left[a_{1} b_{1}+a_{2} b_{2}\right] \omega .
\end{aligned}
$$

Here, we use the fact $1+\alpha+1+\alpha=(1+\alpha)(1+\alpha) \subseteq 1+\alpha$.
Now, suppose true for $n$ and let $a_{1}, a_{2}, \ldots, a_{n}, a_{n+1}, b_{1}, b_{2}, \ldots, b_{n}, b_{n+1} \in F$.

$$
\begin{aligned}
& \left(a_{1}+b_{1} \omega\right)^{2}+\left(a_{2}+b_{2} \omega\right)^{2}+\ldots+\left(a_{n}+b_{n} \omega\right)^{2}+\left(a_{n+1}+b_{n+1} \omega\right)^{2} \\
& =\left[\left(a_{1}+b_{1} \omega\right)^{2}+\left(a_{2}+b_{2} \omega\right)^{2}+\ldots+\left(a_{n}+b_{n} \omega\right)^{2}\right]+\left(a_{n+1}+b_{n+1} \omega\right)^{2} \\
& =(1+\alpha)+2\left[a_{1} b_{1}+\ldots+a_{n} b_{n}\right] \omega+1+2 a_{n+1} b_{n+1} \omega+\alpha \\
& =(1+\alpha)+2\left[a_{1} b_{1}+\ldots+a_{n} b_{n}+a_{n+1} b_{n+1}\right] \omega,
\end{aligned}
$$

as desired.
b-Using the previous item, we get

$$
(1-\omega)^{2}+(1+\omega)^{2}=(1+\alpha)+2[1-1] \omega=(1+\alpha)+F \cdot \omega
$$

Moreover, $(1+\alpha)+F \cdot \omega \subseteq \sum S_{F}(\omega)^{2} \backslash\{0\}$, completing the proof.
c - Let $a, b \in F$. Then

$$
(a+b \omega)^{2}=a^{2}+a b \omega+a b \omega+b^{2} \omega^{2}=(1+\alpha)+2 a b \omega .
$$

Then $F(\omega)^{2} \subseteq(1+\alpha)+(2) F \cdot \omega$. Conversely, let $t \in(1+\alpha)+(2) F \cdot \omega$. Then $t \in(1+\alpha)+2 a \omega$ for some $a \in F$. Since

$$
(1+\alpha)+2 a \omega=(1+a \omega)^{2},
$$

we get $t \in(1+a \omega)^{2} \subseteq F(\omega)^{2}$, completing the proof.
d - Just use items (a), (b) and (c).
e - If $-1 \in F(\omega)^{2}$, then $-1 \in x+2 z \omega$ for some $x \in 1+\alpha$ and $z \in F$. If $z=0$ then $-1=1+\alpha$, contradiction. If $z \neq 0$, then

$$
0 \in 1-1 \subseteq 1+x+2 z \omega \subseteq 1+1+\alpha+2 z \omega .
$$

Then $0 \in y+z \omega$ for some $y \in 1+1+\alpha$, and then, $0 \in \operatorname{ev}(g(X), \omega)$, for $g(X)=y+z X$, contradicting the fact that $f(X)=X^{2}-\alpha=\operatorname{Irr}(F, \omega)$.
f - Just use the same argument of item (d).
g- Let $t \in a+b$ for some $t \in x+2 y \omega$, with $y \neq 0$. Then

$$
-b \in a-t \subseteq a-x-2 y \omega
$$

and then,

$$
0 \in b-b \subseteq b+a-x-2 y \omega .
$$

Therefore exists some $d \in b+a-x$ such that $0 \in d-2 y \omega$, which imply that $0 \in \operatorname{ev}(g(X), \omega)$ for $g(X)=d-z X$ for some $z \in 2 y$ with $z \neq 0$, contradiction. This morphism is not injective because if $a \in 1+\alpha$ then $[a]=[1]$ in $S_{F}(\omega)$.
h- If $a, b \in a+b$ for all $a, b \in \dot{F}$, then (2)F=F. Hence $F(\omega)^{2}=\sum F(\omega)^{2}$, which imply $S_{F}(\omega)$ real reduced $(1+1=\{1\})$.

Remark 6.2.3. It is not an easy task to find the description of $S_{F}(\omega)$ in the language/theory of special groups. Also, it is not clear if such description provides an advantage in terms of comprehension of the following results.

If $\varphi=\left\langle a_{1}, \ldots, a_{n}\right\rangle$ is a form on $F$, we denote the form $[\varphi]$ on $S_{F}(\omega)$ simply by $[\varphi]:=\left\langle\left[a_{1}\right], \ldots,\left[a_{n}\right]\right\rangle$. We say that $[\varphi]$ is the equivalence class of $\varphi$ in $S_{F}(\omega)$. Of course, if $\varphi=\varphi_{1} \oplus \varphi_{2}\left(\right.$ or $\left.\varphi=\varphi_{1} \otimes \varphi_{2}\right)$ then $[\varphi]=\left[\varphi_{1}\right] \oplus\left[\varphi_{2}\right]$ (or $[\varphi]=\left[\varphi_{1}\right] \otimes\left[\varphi_{2}\right]$ ). We have the following useful consequences of Proposition 6.2 .2 (which we will use freely):

## Remark 6.2.4.

$a$ - If $\varphi \equiv \psi$ on $F$ then $[\varphi] \equiv[\psi]$ on $S_{F}(\omega)$;
$b$ - if $\varphi$ is isotropic on $F$ then $[\varphi]$ is isotropic on $S_{F}(\omega)$;
$c-$ if $[\varphi]$ is anisotropic on $S_{F}(\omega)$ then if $\varphi$ is anisotropic on $F$.
Definition 6.2.5 (Rooted Superfield). A superfield $F$ is rooted if

$$
\{a, b\} \subseteq a+b \text { for all } a, b \in F \backslash\{0\} .
$$

Theorem 6.2.6. Let $F$ be a formally real pre-special hyperfield and $\omega \in \bar{F} \backslash F$ be a root of $f(X)=X^{2}-\alpha \in F[X]$. Suppose that $-1 \notin 1+\alpha$. Then $S_{F}(\omega)$ is a formally real pre-special hyperfield. Moreover if $F$ is rooted then $S_{F}(\omega)=S_{F}^{\text {red }}(\omega)$, and in particular, $S_{F}(\omega)$ is a real reduced hyperfield.
Proof. We already know (using Theorem 6.1.7 and item (e) of Proposition 6.2.2) that $S_{F}(\omega)$ is a formally real pre-special hyperfield. If $F$ is rooted, then by item (h) of Proposition 6.2 .2 we have the desired.

Example 6.2.7 (Quadratic Field Extensions and Quadratic Hyperfield Extensions). Let $F$ be a formally real field and $p \in F^{*}$ such that $x^{2}-p$ has no roots in $F$. Consider $K=F(\sqrt{p})$. Of course, we have two special multifields (and special groups) $G(K):=K / m\left(K^{2}\right)^{*}$ and $G_{r e d}(K)=$ $K / m\left(\sum K^{2}\right)^{*}$. Note that if $a+b \sqrt{p} \in \dot{K}$, then

$$
(a+b \sqrt{p})^{2}=a^{2}+p b^{2}+2 a b \sqrt{p} \in D_{F}(\langle 1, p\rangle)+\sqrt{p} \cdot F,
$$

where $D_{F}(\langle 1, p\rangle)$ is the usual set of representatives of the $F$-quadratic form $\langle 1, p\rangle$ :

$$
D_{F}(\langle 1, p\rangle):=\left\{x^{2}+y^{2} p: x, y \in \dot{F}\right\} .
$$

In other words,

$$
K^{2} \backslash\{0\} \subseteq D_{F}(\langle 1, p\rangle)+\sqrt{p} \cdot F .
$$

Moreover,

$$
\begin{aligned}
(a+b \sqrt{p})^{2}+(c+d \sqrt{p})^{2} & =\left(a^{2}+p b^{2}+2 p a b \sqrt{p}\right)+\left(c^{2}+p d^{2}+2 c d \sqrt{p}\right) \\
& =\left(a^{2}+p b^{2}+c^{2}+p d^{2}\right)+2(a b+c d) \sqrt{p} .
\end{aligned}
$$

Using the fact that $D_{F}(\langle 1, p\rangle) \cdot D_{F}(\langle 1, p\rangle) \subseteq D_{F}(\langle 1, p\rangle)$ and for $a, b, c, d \neq 0$,

$$
a^{2}+p b^{2}+c^{2}+p d^{2} \in D_{F}(\langle 1, p, 1, p\rangle)=D_{F}(\langle 1, p\rangle \otimes\langle 1, p\rangle)=D_{F}(\langle 1, p\rangle) \cdot D_{F}(\langle 1, p\rangle),
$$

we conclude by induction (and a case analysis involving $0 \in\{a, b, c, d\}$ ) that

$$
\sum K^{2} \backslash\{0\} \subseteq D_{F}(\langle 1, p\rangle)+\sqrt{p} \cdot F .
$$

So, let $Q_{p}:=D_{K}(\langle 1, p\rangle)+\sqrt{p} \cdot F$. Then $Q_{p} \cdot Q_{p}$ is a multiplicative set containing $\sum K^{2}$. Define $G_{\sqrt{p}}(K):=K /{ }_{m} Q_{p}$. Then $G_{\sqrt{p}}(K)$ is a reduced special group such that

$$
\begin{equation*}
G(K) \rightarrow G_{r e d}(K) \rightarrow G_{\sqrt{p}}(K) . \tag{}
\end{equation*}
$$

Moreover,

$$
G_{\sqrt{p}}(K) \cong S_{K / m\left(K^{2}\right)^{*}}(\sqrt{p}),
$$

i.e, the hyperfield of Theorem 6.2.6. We say that $K$ is $p$-special if $G_{\sqrt{p}}(K) \cong G_{r e d}(K)$.
exquad2

Example 6.2.8 (The Special Group of a quadratic extension). Let $F$ be a formally real field and $p \in F^{*}$ such that $x^{2}-p$ has no roots in $F$. Consider $K=F(\sqrt{p})$. Using the calculations of Example 6.2.7 we have

$$
\begin{aligned}
\dot{K}^{2} & =D_{F}(1, p)+\left\{x^{2}+y b^{2}+z \sqrt{p}: x, y \neq \dot{F} \text { and } z=(x+y)^{2}-\left(x^{2}+y^{2}\right)\right\} \\
& =\left\{x^{2}+y b^{2}+z \sqrt{p}: x, y \neq F \text { are not both } 0 \text { and } z=(x+y)^{2}-\left(x^{2}+y^{2}\right)\right\}
\end{aligned}
$$

In this sense, for $a, b, c, d \in \dot{K}$, what means $\langle a, b\rangle \equiv_{K}\langle c, d\rangle$ in terms of the isometry relation on $F$ ?

By Lemma 1.5(a) of [28] we have

$$
\langle a, b\rangle \equiv_{K}\langle c, d\rangle \text { iff } a b=c d \text { and } a c \in D_{K}(1, c d) .
$$

Lets first understand what means $\beta \in D_{K}(1, \alpha)$ for $\alpha, \beta \in \dot{K}$. By definition,

$$
\beta \in D_{K}(1, \alpha) \text { iff } \beta=x^{2}+\alpha y^{2}, x, y \in \dot{K}
$$

Write $\alpha=\alpha_{1}+\alpha_{2} \sqrt{p}, \beta=\beta_{1}+\beta_{2} \sqrt{p}, x=x_{1}+x_{2} \sqrt{p}$ and $y=y_{1}+y_{2} \sqrt{p}$. Then

$$
\begin{aligned}
\beta & =x^{2}+\alpha y^{2} \Leftrightarrow \\
\beta_{1}+\beta_{2} \sqrt{p} & =\left(x_{1}+x_{2} \sqrt{p}\right)^{2}+\left(\alpha_{1}+\alpha_{2} \sqrt{p}\right)\left(y_{1}+y_{2} \sqrt{p}\right)^{2} \Leftrightarrow \\
\beta_{1}+\beta_{2} \sqrt{p} & =\left(x_{1}^{2}+p x_{2}^{2}+2 x_{1} x_{2} \sqrt{p}\right)+\left(\alpha_{1}+\alpha_{2} \sqrt{p}\right)\left(y_{1}^{2}+p y_{2}^{2}+2 y_{1} y_{2} \sqrt{p}\right) \Leftrightarrow \\
\beta_{1}+\beta_{2} \sqrt{p} & =\left(x_{1}^{2}+p x_{2}^{2}+\alpha_{1} y_{1}^{2}+\alpha_{1} p y_{2}^{2}+2 p \alpha_{2} y_{1} y_{2}\right)+\left(2 x_{1} x_{2}+2 \alpha_{1} y_{1} y_{2}+\alpha_{2} y_{1}^{2}+\alpha_{2} p y_{2}^{2}\right) \sqrt{p} \\
& \Leftrightarrow\left\{\begin{array}{l}
\beta_{1}=x_{1}^{2}+p x_{2}^{2}+\alpha_{1} y_{1}^{2}+\alpha_{1} p y_{2}^{2}+2 p \alpha_{2} y_{1} y_{2} \\
\beta_{2}=2 x_{1} x_{2}+2 \alpha_{1} y_{1} y_{2}+\alpha_{2} y_{1}^{2}+\alpha_{2} p y_{2}^{2}
\end{array}\right. \\
& \Leftrightarrow\left\{\begin{array}{l}
\beta_{1}+\beta_{2}=\left(x_{1}+x_{2}\right)^{2}+\left(\alpha_{1}+\alpha_{2} p\right)\left(y_{1}+y_{2}\right)^{2}+(p-1)\left(x_{2}^{2}+\alpha_{1} y_{2}^{2}-\alpha_{2} y_{1}^{2}\right) \\
\beta_{1}-\beta_{2}=\left(x_{1}-x_{2}\right)^{2}+\left(\alpha_{1}-\alpha_{2} p\right)\left(y_{1}-y_{2}\right)^{2}+(p-1)\left(x_{2}^{2}+\alpha_{1} y_{2}^{2}+\alpha_{2} y_{1}^{2}\right)
\end{array}\right.
\end{aligned}
$$

Then

$$
\beta=x^{2}+a y^{2} \Leftrightarrow\left\{\begin{array}{l}
\beta_{1}+\beta_{2}=\left(x_{1}+x_{2}\right)^{2}+\left(a_{1}+a_{2} p\right)\left(y_{1}+y_{2}\right)^{2}+(p-1)\left(x_{2}^{2}+a_{1} y_{2}^{2}-a_{2} y_{1}^{2}\right)  \tag{eqkrep 1}\\
\beta_{1}-\beta_{2}=\left(x_{1}-x_{2}\right)^{2}+\left(a_{1}-a_{2} p\right)\left(y_{1}-y_{2}\right)^{2}+(p-1)\left(x_{2}^{2}+a_{1} y_{2}^{2}+a_{2} y_{1}^{2}\right)
\end{array}\right.
$$

For the discriminant part, let $a, b, c, d \in \dot{K}$ with $a=a_{1}+a_{2} \sqrt{p}, b=b_{1}+b_{2} \sqrt{p}, c=c_{1}+c_{2} \sqrt{p}$
and $d=d_{1}+d_{2} \sqrt{p}$ for suitable $a_{i}, b_{i}, c_{i}, d_{i} \in F(i=1,2)$. We have

$$
\begin{aligned}
a b=c d & \Leftrightarrow\left(a_{1}+a_{2} \sqrt{p}\right)\left(b_{1}+b_{2} \sqrt{p}\right)=\left(c_{1}+c_{2} \sqrt{p}\right)\left(d_{1}+d_{2} \sqrt{p}\right) \\
& \Leftrightarrow\left(a_{1} b_{1}+p a_{2} b_{2}\right)+\left(a_{1} b_{2}+a_{2} b_{1}\right) \sqrt{p}=\left(c_{1} d_{1}+p c_{2} d_{2}\right)+\left(c_{1} d_{2}+c_{2} d_{1}\right) \sqrt{p} \\
& \Leftrightarrow\left\{\begin{array}{l}
a_{1} b_{1}+p a_{2} b_{2}=c_{1} d_{1}+p c_{2} d_{2} \\
a_{1} b_{2}+a_{2} b_{1}=c_{1} d_{2}+c_{2} d_{1}
\end{array}\right. \\
& \Leftrightarrow\left\{\begin{array}{l}
\left(a_{1} b_{1}+p a_{2} b_{2}\right)+\left(a_{1} b_{2}+a_{2} b_{1}\right)=\left(c_{1} d_{1}+p c_{2} d_{2}\right)+\left(c_{1} d_{2}+c_{2} d_{1}\right) \\
\left(a_{1} b_{1}+p a_{2} b_{2}\right)-\left(a_{1} b_{2}+a_{2} b_{1}\right)=\left(c_{1} d_{1}+p c_{2} d_{2}\right)-\left(c_{1} d_{2}+c_{2} d_{1}\right)
\end{array}\right. \\
& \Leftrightarrow\left\{\begin{array}{l}
\left(a_{1}+a_{2}\right)\left(b_{1}+b_{2}\right)+(p-1) a_{2} b_{2}=\left(c_{1}+c_{2}\right)\left(d_{1}+d_{2}\right)+(p-1) c_{2} d_{2} \\
\left(a_{1}-a_{2}\right)\left(b_{1}-b_{2}\right)+(p-1) a_{2} b_{2}=\left(c_{1}-c_{2}\right)\left(d_{1}-d_{2}\right)+(p-1) c_{2} d_{2}
\end{array}\right.
\end{aligned}
$$

Then

$$
a b=c d \Leftrightarrow\left\{\begin{array}{l}
\left(a_{1}+a_{2}\right)\left(b_{1}+b_{2}\right)+(p-1) a_{2} b_{2}=\left(c_{1}+c_{2}\right)\left(d_{1}+d_{2}\right)+(p-1) c_{2} d_{2}  \tag{}\\
\left(a_{1}-a_{2}\right)\left(b_{1}-b_{2}\right)+(p-1) a_{2} b_{2}=\left(c_{1}-c_{2}\right)\left(d_{1}-d_{2}\right)+(p-1) c_{2} d_{2}
\end{array}\right.
$$

Now, since ac $=\left(a_{1} c_{1}+p a_{2} c_{2}\right)+\left(a_{1} c_{2}+a_{2} c_{1}\right) \sqrt{p}$ and ad $=\left(a_{1} d_{1}+p a_{2} d_{2}\right)+\left(a_{1} c_{2}+a_{2} d_{1}\right) \sqrt{p}$, using Equations 6.1 and 6.2 (with $\beta=a c$ and $\alpha=a d$ in Equation 6.1), we have the following characterization:
$\langle a, b\rangle \equiv_{K}\langle c, d\rangle$ if and only if there exists $x_{1}, x_{2}, y_{1}, y_{2} \in F$ such that $\left(x_{1}, x_{2}\right),\left(y_{1}, y_{2}\right) \neq(0,0)$ and

$$
\left\{\begin{aligned}
\left(a_{1}+a_{2}\right)\left(b_{1}+b_{2}\right)+(p-1) a_{2} b_{2} & =\left(c_{1}+c_{2}\right)\left(d_{1}+d_{2}\right)+(p-1) c_{2} d_{2} \\
\left(a_{1}-a_{2}\right)\left(b_{1}-b_{2}\right)+(p-1) a_{2} b_{2} & =\left(c_{1}-c_{2}\right)\left(d_{1}-d_{2}\right)+(p-1) c_{2} d_{2} \\
\left(a_{1}+a_{2}\right)\left(c_{1}+c_{2}\right)+(p-1) a_{2} c_{2} & =\left(x_{1}+x_{2}\right)^{2}+\left[p\left(a_{1}+a_{2}\right)\left(d_{1}+d_{2}\right)-(p-1) a_{1} d_{1}\right]\left(y_{1}+y_{2}\right)^{2} \\
& +(p-1)\left[x_{2}^{2}+\left(a_{1} d_{1}+p a_{2} d_{2}\right) y_{2}^{2}-\left(a_{1} c_{2}+a_{2} d_{1}\right) y_{1}^{2}\right] \\
\left(a_{1}-a_{2}\right)\left(c_{1}-c_{2}\right)+(p-1) a_{2} c_{2} & =\left(x_{1}-x_{2}\right)^{2}+\left[p\left(a_{1}-a_{2}\right)\left(d_{1}-d_{2}\right)-(p-1) a_{1} d_{1}\right]\left(y_{1}-y_{2}\right)^{2} \\
& +(p-1)\left[x_{2}^{2}+\left(a_{1} d_{1}+p a_{2} d_{2}\right) y_{2}^{2}+\left(a_{1} c_{2}+a_{2} d_{1}\right) y_{1}^{2}\right]
\end{aligned}\right.
$$

Manipulating Equation 6.2 we get a very similar system to describe when $\bar{a}=\bar{b}$ in $K / m \dot{K}^{2}, a, b \in K$.
In the sequel, we want to iterate de construction $S_{F}(\omega)$. Let $\alpha, \beta \in \dot{F} \backslash\{-1\}$. The properties of Proposition 6.2.2 are valid if we change $F$ by $S_{F}(\sqrt{\alpha})(\sqrt{\beta})$ (or $S_{F}(\sqrt{\beta})(\sqrt{\alpha})$ ).

Theorem 6.2.9. Let $F$ be a special hyperfield and $\alpha, \beta \in \dot{F} \backslash\{ \pm 1\}$. Then

$$
S_{S_{F}(\sqrt{\alpha})}(\sqrt{\beta}) \cong S_{S_{F}(\sqrt{\beta})}(\sqrt{\alpha}) .
$$

Proof. We already know that

$$
F(\sqrt{\alpha})(\sqrt{\beta}) \cong F(\sqrt{\beta})(\sqrt{\alpha})
$$

and

$$
F(\sqrt{\alpha})(\sqrt{\beta})=F+F \sqrt{\alpha}+F \sqrt{\beta}+F \sqrt{\alpha} \sqrt{\beta} .
$$

Let $\varphi: F(\sqrt{\alpha})(\sqrt{\beta}) \rightarrow F(\sqrt{\beta})(\sqrt{\alpha})$ be an isomorphism and denote $q_{1}: F(\sqrt{\alpha})(\sqrt{\beta}) \rightarrow S_{F}(\sqrt{\alpha})(\sqrt{\beta})$ and $q_{2}: F(\sqrt{\beta})(\sqrt{\alpha}) \rightarrow S_{F}(\sqrt{\beta})(\sqrt{\alpha})$ the natural projections. For instance, $q_{1}$ is given by the rule

$$
q_{1}\left(a_{0}+a_{1} \sqrt{\alpha}+a_{2} \sqrt{\beta}+a_{3} \sqrt{\alpha} \sqrt{\beta}\right):=\left[a_{0}+a_{1} \sqrt{\alpha}\right]+\left[a_{2}+a_{3} \sqrt{\alpha}\right] \sqrt{\beta} \in S_{F}(\sqrt{\alpha})(\sqrt{\beta}) .
$$

Similarly for $q_{2}$.
Now, let $\pi_{1}: S_{F}(\sqrt{\alpha})(\sqrt{\beta}) \rightarrow S_{S_{F}(\sqrt{\alpha})}(\sqrt{\beta})$ and $\pi_{2}: S_{F}(\sqrt{\beta})(\sqrt{\alpha}) \rightarrow S_{S_{F}(\sqrt{\beta})}(\sqrt{\alpha})$ be the quotient morphisms. Since $\varphi\left[F(\sqrt{\alpha})(\sqrt{\beta})^{2}\right]=F(\sqrt{\beta})(\sqrt{\alpha})^{2}$ and $q_{2}\left[F(\sqrt{\beta})(\sqrt{\alpha})^{2}\right] \subseteq\left(S_{F}(\sqrt{\beta})(\sqrt{\alpha})\right)^{2}$, we have that $\pi_{2} \circ q_{2} \circ \varphi: F(\sqrt{\alpha})(\sqrt{\beta}) \rightarrow S_{S_{F}(\sqrt{\beta})}(\sqrt{\alpha})$ is a morphism such that

$$
\pi_{2} \circ q_{2} \circ \varphi[F(\sqrt{\alpha})(\sqrt{\beta})]=\{1\} .
$$

By the universal property there is an unique morphism $\varphi_{\alpha \beta}: S_{S_{F}(\sqrt{\alpha})}(\sqrt{\beta}) \rightarrow S_{S_{F}(\sqrt{\beta})}(\sqrt{\alpha})$. Using the same argument, there is an unique morphism $\varphi_{\beta \alpha}: S_{S_{F}(\sqrt{\beta})}(\sqrt{\alpha}) \rightarrow S_{S_{F}(\sqrt{\alpha})}(\sqrt{\beta})$. The universal property forces $\varphi_{\alpha \beta} \circ \varphi_{\beta \alpha}=i d$ and $\varphi_{\beta \alpha} \circ \varphi_{\alpha \beta}=i d$.

With Theorem 6.2.9 we are able to properly iterate the construction $S_{F}(\omega)$. For $a_{1}, \ldots, a_{n} \in \dot{F}$, we define recursively:

$$
\begin{aligned}
S_{F\left(\sqrt{a_{1}}, \sqrt{a_{2}}\right)} & :=S_{S_{F}\left(\sqrt{a_{1}}\right)}\left(\sqrt{a_{2}}\right) \\
S_{F\left(\sqrt{a_{1}}, \ldots, \sqrt{a_{n+1}}\right)} & :=S_{S_{F}\left(\sqrt{a_{1}}, \ldots, \sqrt{a_{n}}\right)}\left(\sqrt{a_{n+1}}\right) .
\end{aligned}
$$

Corollary 6.2.10. Let $F$ be a special hyperfield, $\alpha_{1}, \alpha_{2}, \ldots, \alpha_{n} \in \dot{F}$, and $\sigma \in S_{n}$. Then

$$
S_{F\left(\sqrt{a_{1}}, \ldots, \sqrt{a_{n}}\right)} \cong S_{F\left(\sqrt{a_{\sigma(1)}}, \ldots, \sqrt{a_{\sigma(n)}}\right)} .
$$

It is important to comprehend the distinction between $S_{F\left(\sqrt{a_{1}}, \ldots, \sqrt{a_{n}}\right)}\left(\sqrt{a_{n+1}}\right)$ and $S_{F\left(\sqrt{a_{1}}, \ldots, \sqrt{a_{n+1}}\right)}$ : $S_{F\left(\sqrt{a_{1}}, \ldots, \sqrt{a_{n}}\right)}\left(\sqrt{a_{n+1}}\right)$ is an algebraic extension of $S_{F\left(\sqrt{a_{1}}, \ldots, \sqrt{a_{n}}\right)}$ from which $\sqrt{a_{n+1}}$ is algebraic. On the other hand,

$$
S_{F\left(\sqrt{a_{1}}, \ldots, \sqrt{a_{n+1}}\right)}:=S_{F\left(\sqrt{a_{1}}, \ldots, \sqrt{a_{n}}\right)}\left(\sqrt{a_{n+1}}\right) / m\left(\left(S_{F\left(\sqrt{a_{1}}, \ldots, \sqrt{a_{n}}\right)}\left(\sqrt{a_{n+1}}\right)\right)^{2} \backslash\{0\}\right)
$$

We want to describe the isometry relation $\equiv_{S_{F}(\omega)}$ in $S_{F}(\omega)$ in terms of isometry relation $\equiv_{F}$ in $F$. We begin this investigation with some general results.
iso0
Theorem 6.2.11. Let $a, b \in F$. Then $[a]=[b]$ in $S_{F}(\omega)$ iff there is $s, t \in 1+\alpha$ with as $=b t$ (or $a=b s t$ ).

Proof. $(\Rightarrow)$ Suppose that $[a]=[b]$ in $S_{F}(\omega)$. Then there exist $X, Y \subseteq F(\omega)^{2} \backslash\{0\}$ with $a X=b Y$. Let $r_{1}+s_{1} \omega \in X$, with $r \in 1+\alpha$. Then

$$
a\left(r_{1}+s_{1} \omega\right)=a r_{1}+a s_{1} \omega \in b Y,
$$

and there exist $r_{2}+s_{2} \omega \in Y$ with

$$
\emptyset \neq\left(a r_{1}+a s_{1} \omega\right) \cap\left(b\left(r_{2}+s_{2} \omega\right)\right) .
$$

But $b\left(r_{2}+s_{2} \omega\right)=\left\{b r_{2}+b s_{2} \omega\right\}$. Then $a r_{1}+a s_{1} \omega=b r_{2}+b s_{2} \omega$, which imply $a r_{1}=b r_{2}$ and $a s_{1}=b s_{2}$.
$(\Leftarrow)$ Immediate since $1+\alpha \subseteq F(\omega)^{2} \backslash\{0\}$.
Theorem 6.2.12. Let $a, b \in F$. Then $[a]=[b]$ in $S_{F\left(\sqrt{a_{1}}, \ldots, \sqrt{a_{n}}\right)}$ iff there is $s, t \in D_{F}\left(\left\langle\left\langle\alpha_{1}, \alpha_{2}, \ldots, \alpha_{n}\right\rangle\right\rangle\right)$ 1 such that $a s=$ bt (or $a=$ bst or even $a b \in D_{F}\left(\left\langle\left\langle\alpha_{1}, \alpha_{2}, \ldots, \alpha_{n}\right\rangle\right\rangle\right)$ ).

[^14]Proof. Since $S_{F\left(\sqrt{a_{1}}, \ldots, \sqrt{a_{n}}\right)}$ is a real reduced hyperfield and in $S_{F\left(\sqrt{a_{1}}, \ldots, \sqrt{a_{n}}\right)},\left\{\left[\alpha_{1}\right], \ldots,\left[\alpha_{n}\right]\right\}=\{[1]\}$ we only need to prove $(\Rightarrow)$.

Let $a, b \in F$. If $[a]=[b]$ in $S_{F\left(\sqrt{\alpha_{1}}, \sqrt{\alpha_{2}}\right)}=S_{S_{F\left(\sqrt{\alpha_{1}}\right)}}\left(\sqrt{\alpha_{2}}\right)$, then by Theorem 6.2.11 (changing $F$ by $S_{F\left(\sqrt{\alpha_{1}}\right)}$ ) we have $[a b] \in[1]+\left[\alpha_{2}\right]$ (in $S_{F\left(\sqrt{\alpha_{1}}\right)}$ ). Then there exist $s \in 1+\alpha_{1}$ such that

$$
a b s \in\left(1+\alpha_{1}\right)+\alpha_{2}\left(1+\alpha_{1}\right)=D_{F}\left(\left\langle 1, \alpha_{1}, \alpha_{2}, \alpha_{1} \alpha_{2}\right\rangle\right)=D_{F}\left(\left\langle\left\langle\alpha_{1}, \alpha_{2}\right\rangle\right\rangle\right) .
$$

This means abs $\in D_{F}\left(\left\langle\left\langle\alpha_{1}, \alpha_{2}\right\rangle\right\rangle\right)$.
Now suppose the desired valid for $n$. By induction hypothesis $[a]=[b]$ in

$$
\left.S_{F\left(\sqrt{a_{1}}, \ldots, \sqrt{a_{n+1}}\right)} \cong S_{\left(S_{F}\left(\sqrt{a_{1}}\right)\right)\left(\sqrt{a_{2}}, \ldots, \sqrt{a_{n+1}}\right)}\right]^{2}
$$

iff

$$
[a b] \in D_{S_{F}\left(\sqrt{a_{1}}\right)}\left(\left\langle\left\langle\left[\alpha_{2}\right],\left[\alpha_{2}\right], \ldots,\left[\alpha_{n}\right]\right\rangle\right\rangle\right) .
$$

This imply

$$
a b \in D_{F}\left(\left\langle\left\langle\alpha_{2}, \alpha_{2}, \ldots, \alpha_{n}\right\rangle\right\rangle\right) \cdot\left(1+\alpha_{1}\right) \subseteq D_{F}\left(\left\langle\left\langle\alpha_{1}, \alpha_{2}, \ldots, \alpha_{n+1}\right\rangle\right\rangle\right) .
$$

cor 33
Corollary 6.2.13. Let $F$ be a special hyperfield and $\alpha_{1}, \alpha_{2}, \ldots, \alpha_{n} \in \dot{F} \backslash\{ \pm 1\}$. Then $S_{F\left(\sqrt{a_{1}}, \ldots, \sqrt{a_{n}}\right)}$ is formally real iff $-1 \notin D_{F}\left(\left\langle\left\langle\alpha_{1}, \alpha_{2}, \ldots, \alpha_{n}\right\rangle\right\rangle\right)$.

Proof. Since $S_{F\left(\sqrt{a_{1}}, \ldots, \sqrt{a_{n}}\right)}$ is a real reduced hyperfield, we have $S_{F\left(\sqrt{a_{1}}, \ldots, \sqrt{a_{n}}\right)}$ formally real iff $[1] \neq[-1]$, which by Theorem 6.2.12 occurs iff $-1 \notin D\left(\left\langle\left\langle\alpha_{1}, \alpha_{2}, \ldots, \alpha_{n}\right\rangle\right\rangle\right)$.
iso2
Theorem 6.2.14. Let $a, b, c, d \in \dot{F}$. Then $\langle[a],[b]\rangle \equiv_{S_{F}(\omega)}\langle[c],[d]\rangle$ iff $\langle a r, b s\rangle \equiv_{F}\langle c, d t\rangle$ for some $r, s, t \in 1+\alpha$.

Proof. $(\Rightarrow)$ Let $\langle[a],[b]\rangle \equiv_{S_{F}(\omega)}\langle[c],[d]\rangle$. Then $[a][b]=[c][d]$ and $[a][c] \in 1+[c][d]$. Hence, there are $v, w \in 1+\alpha$ and $x \in S=F(\omega)^{2} \backslash\{0\}$ with $a b v=c d w$ (or $a b v w=c d$ ) and $a c x \in S+c d S$. Write $x=x_{1}+x_{2} \omega$. We have

$$
a v w c x_{1} x \in v x_{1}(S+c d S) \subseteq v x_{1} S+c d v x_{1} S \subseteq S+c d S .
$$

Then

$$
\operatorname{avwcx}_{1}\left(x_{1}+x_{2} \omega\right) \in\left(y_{1}+y_{2} \omega\right)+c d\left(z_{1}+z_{2} \omega\right)
$$

for some $y_{1}, z_{1} \in 1+\alpha$ and $y_{2}, z_{2} \in F$. This means

$$
a v w c+a v w c x_{1} x_{2} \omega \in\left(y_{1}+c d z_{1}\right)+\left(y_{2}+c d z_{2}\right) \omega
$$

and then, $a v w c \in y_{1}+c d z_{1}$ and $a v w c x_{1} x_{2} \in y_{2}+c d z_{2}$. Then $a v w c y_{1} \in 1+c d y_{1} z_{1}$ or equivalently, $\left(a v w y_{1}\right) c \in 1+c\left(d y_{1} z_{1}\right)$. Therefore $\left(a v w y_{1}\right)\left(b z_{1}\right)=c\left(d y_{1} z_{1}\right)$ with $\left(a v w y_{1}\right) c \in 1+c\left(d y_{1} z_{1}\right)$, which means $\left\langle a v w y_{1}, b z_{1}\right\rangle \equiv_{F}\left\langle c, d y_{1} z_{1}\right\rangle$. Putting $r=v w y_{1}, s=z_{1}$ and $t=y_{1} z_{1}$ we get the desired.
$(\Leftarrow)$ Immediate.
Then for all $a, b, c, d \in \dot{F}$ are equivalent:

$$
\begin{aligned}
\text { i- }\langle[a],[b]\rangle & \equiv_{S_{F}(\omega)}\langle[c],[d]\rangle ; \\
\text { ii- }\langle a r, b s\rangle & \equiv_{F}\langle c t, d\rangle \text { for some } r, s, t \in 1+\alpha .
\end{aligned}
$$

[^15]iii- $\langle a r, b s\rangle \equiv_{F}\langle c, d t\rangle$ for some $r, s, t \in 1+\alpha$.
iv- $\langle a, b r\rangle \equiv_{F}\langle c s, d t\rangle$ for some $r, s, t \in 1+\alpha$.
$\mathrm{v}-\langle a r, b\rangle \equiv_{F}\langle c s, d t\rangle$ for some $r, s, t \in 1+\alpha$.

### 6.3 Expanding the Arason-Pfister Hauptsatz and consequences

Hauptsatz-section
As an application of the former developed results, we reserve this section to expand one of the more emblematic questions/results in algebraic theory of quadratic forms: the so-called Arason Pfister Hauptsatz (APH).

In the sequel we present a brief historic of APH.
In [52], a 1970 paper of John Milnor seminal to the algebraic theory of quadratic forms over fields, the author poses two questions concerning the class of fields of characteristic $\neq 2$ (positively solved in the paper in many instances). One of the question was concerning the so called "Milnor's conjectures for the graded cohomology ring and for the graded Witt ring" that Voevodsky et al. solved around 2000. The other question asked if for every such field $F$, the intersection $\bigcap_{n \in \mathbb{N}} I^{n}(F)$ contains only $0 \in W(F)$, where $I^{n}(F)$ is the n-th power of the fundamental ideal $I(F)$ of the Witt ring of $F(I(F)=\{$ even dimensional anisotropic forms over $F\}$ ).

In the subsequent year, J. Arason and A. Pfister solved this question as an immediate corollary of the nowadays called "Arason-Pfister Hauptsatz" (APH):
([7]) Let $\phi \neq \emptyset$ be an anisotropic form. If $\phi \in I^{n}(F)$, then $\operatorname{dim}(\phi) \geq 2^{n}$.
The theory of special groups, an abstract (first-order) theory of quadratic forms developed by Dickmann-Miraglia since the middle of the 1990s, allows a functorial encoding of the algebraic theory of quadratic forms of fields (with char $\neq 2$ ). In [28], Dickmann-Miraglia, restated the APH to the setting of special groups and, employing boolean theoretic methods to define and calculate the Stiefel-Whitney and the Horn-Tarski invariants of a special group, establish a generalization of the APH to the setting of reduced special groups, in particular proving a different proof of the APH for formally real pythagorian fields.

Now, as an application of the previous developed constructions of quadratic extensions in the category of hyperfields, we expand the validity of the Arason-Pfister Hauptsatz for all special hyperfields.

We start establishing the following:

## Notations:

- Let $G$ be a special group and $F=G \dot{\cup}\{0\}$ its special hyperfield associated. In particular, if $\alpha \in F, \alpha \neq 0,1$, the polynomial $f(X) \in F[X], f(X)=X^{2}-\alpha$ has no roots in $F$ (basically because $\alpha^{2}=1$ for all $\left.\alpha \in F^{*}\right)$.
- Let $\varphi, \psi$ be forms on a special hyperfield $F$. We say that $\varphi$ and $\psi$ are Witt equivalent, notation $\varphi \approx_{W, F} \psi$ iff there exist non negative integers $k, l$ such that $k\langle 1,-1\rangle \oplus \varphi \equiv_{F} l\langle 1,-1\rangle \oplus \psi$. By Witt's Decomposition, if $\varphi$ is a form on $F$, there are unique forms $\varphi_{a n}, \varphi_{h i p}, \varphi_{0}$ (up to isometry) with $\varphi \equiv \varphi_{a n} \oplus \varphi_{h i p} \oplus \varphi_{0}, \varphi_{a n}$ anisotropic, $\varphi_{h i p}$ hyperbolic and $\varphi_{0}$ totally isotropic. We define $\operatorname{dim}_{W, F}(\varphi):=\operatorname{dim}\left(\varphi_{a n}\right)$.
- Let $F$ be a special hyperfield. For each $n \in \mathbb{N}$ consider the statement:
$A P_{F}(n)$ For each $\varphi=\left\langle a_{1}, \cdots, a_{k}\right\rangle$, a non-empty $(k \geq 1)$, regular $\left(a_{i} \in \dot{F}\right)$ and anisotropic form, if $\varphi \in I^{n}(F)$, then $\operatorname{dim}(\varphi) \geq 2^{n}$.


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Remark 6.3.1. Let $G$ be an special group. Recall that (see [28]):
a-A Pfister form of degree $n \geq 1$, with coefficients $a_{1}, \cdots, a_{n} \in G$ is $\left\langle\left\langle a_{1}, \cdots, a_{n}\right\rangle\right\rangle=\otimes_{i=1}^{n}\left\langle 1, a_{i}\right\rangle$.
$b$ - Let $\psi$ be a Pfister form. Then $\psi$ is hyperbolic iff it is isotropic. Moreover, if $G$ is reduced and $-1 \in D_{G}(\psi)$, then $\psi$ is hyperbolic.
$c-I^{n}(G) \subseteq W(G)$ is additively generated by the Pfister forms of degree $n$.
$d$ - If $\varphi \in I^{n}(G) \backslash\{\emptyset\}$, then $\varphi=\varepsilon_{1} \varphi_{1}+\ldots+\varepsilon_{r} \varphi_{r}$, where $r \geq 1$ and $\varepsilon_{j}= \pm 1$ for all $j=1, \ldots, r$. Moreover, if $\varphi$ is anisotropic, we suppose without loss of generality that $\varepsilon_{j}=1$ for all $j=1, \ldots, r$. haup
Theorem 6.3.2 (Arason-Pfister Hauptsatz). Let $F$ be a special hyperfield, then it holds $A P_{F}(n)$, for all $n \geq 0$. In more details: for each $n \geq 0$ and For each $\varphi=\left\langle a_{1}, \cdots, a_{k}\right\rangle$, a non-empty $(k \geq 1)$, regular $\left(a_{i} \in \dot{F}\right)$ and anisotropic form, if $\varphi \in I^{n}(F)$, then $\operatorname{dim}(\varphi) \geq 2^{n} \varphi \in I^{n}(F)$, if $\varphi \neq \emptyset$ is anisotropic, then $\operatorname{dim}_{W, F}(\varphi) \geq 2^{n}$.

Proof. An equivalent way to state this result is the following: if a form $q$ belongs to $I^{n} F$ and $\operatorname{dim} q<2^{n}$ then $q$ must be a hyperbolic form.

Since $\varphi$ is an anisotropic form such that $\varphi \in I^{n} F \backslash\{\emptyset\}$ and $I^{n} F$ is additively generated by the Pfister forms, then there exists $r \geq 1$ and Pfister forms of degree $n, \varphi_{1}, \cdots \varphi_{r}$ such that $\varphi= \pm\left(\varphi_{1}+\ldots+\varphi_{r}\right)$.

Since $\varphi$ is anisotropic, we can suppose without loss of generality that $\varphi=\varphi_{1}+\ldots+\varphi_{r}$ and proceed by induction on $r$.

If $r=1$, then $\varphi=\varphi_{1}$, with $\operatorname{dim}(\varphi)=\operatorname{dim}\left(\varphi_{1}\right)=2^{n}$.
Let $r \geq 2$. If $\varphi_{j}$ is isotropic for all $j=1, \ldots, r$ then $\varphi$ is isotropic (hyperbolic, in fact): this fallows from Witt's cancellation law since $\varphi \oplus k\langle 1,-1\rangle \equiv\left(r 2^{n-1}+m\right)\langle 1,-1\rangle$. So we can suppose without loss of generality that $\varphi_{1}=\left\langle\left\langle a_{1}, \ldots, a_{n}\right\rangle\right\rangle$ is anisotropic.

Suppose $-1 \notin D_{F}\left(\varphi_{1}\right)$. By Corollary 6.2 .13 Let $S_{F}\left(\varphi_{1}\right):=S_{F\left(\sqrt{a_{1}}, \ldots, \sqrt{a_{n}}\right)}$. Then equivalence class of $\varphi$ on $S_{F}\left(\varphi_{1}\right)$ is

$$
[\varphi]=\left[\varphi_{1}+\ldots+\varphi_{r}\right]=\left[\varphi_{1}\right]+\ldots+\left[\varphi_{r}\right]=2^{n}\langle 1\rangle+\left[\varphi_{2}\right]+\ldots+\left[\varphi_{r}\right] .
$$

We already know that $\operatorname{dim}_{W, F}(\varphi) \geq \operatorname{dim}_{W, S_{F}\left(\varphi_{1}\right)}[\varphi]$. Then we have three cases:
I- [ $\varphi$ ] is hyperbolic. Then $\left(\left[\varphi_{2}\right]+\ldots+\left[\varphi_{r}\right]\right)_{a n} \equiv_{S_{F}\left(\varphi_{1}\right)} 2^{n}\langle-1\rangle$. Then
$\operatorname{dim}_{W, F}(\varphi) \geq \operatorname{dim}_{W, F}\left(\varphi_{2}+\ldots+\varphi_{r}\right) \geq \operatorname{dim}_{W, F}\left(\varphi_{2}+\ldots+\varphi_{r}\right)_{a n} \geq \operatorname{dim}_{W, F}\left(\left[\varphi_{2}\right]+\ldots+\left[\varphi_{r}\right]\right)_{a n}=2^{n}$.
II - [ $\varphi$ ] is not hyperbolic and $\left[\varphi_{2}\right]+\ldots+\left[\varphi_{r}\right]$ is anisotropic. Then $\varphi_{2}+\ldots+\varphi_{r}$ is anisotropic. By induction hypothesis we have $\operatorname{dim}_{W, F}\left(\varphi_{2}+\ldots+\varphi_{r}\right) \geq 2^{n}$. Then

$$
\operatorname{dim}_{W, F}(\varphi) \geq \operatorname{dim}_{W, F}\left(\varphi_{2}+\ldots+\varphi_{r}\right) \geq 2^{n}
$$

III - $[\varphi]$ is not hyperbolic and $\left[\varphi_{2}\right]+\ldots+\left[\varphi_{r}\right]$ is isotropic.
Since $[\varphi]$ is not hyperbolic, we can assume that $\left[\varphi_{2}\right]$ is anisotropic (otherwise, if $\left[\varphi_{j}\right]$ is isotropic for all $j=2, \ldots, r$ then $[\varphi]$ is an isotropic Pfister form and then, is also hyperbolic). Denote $F_{1}:=S_{F}\left(\varphi_{1}\right)$. In $S_{F_{1}}\left(\left[\varphi_{2}\right]\right)$ (which is a special hyperfield) look at

$$
\psi_{2}:=\left[\left[\varphi_{2}\right]+\ldots+\left[\varphi_{r}\right]\right]=2^{n}\langle 1\rangle+\left[\varphi_{3}\right]+\ldots+\left[\varphi_{r}\right] \in S_{F_{1}}\left(\left[\varphi_{2}\right]\right) .
$$

For $\psi_{2} \in I^{n}\left(S_{F_{1}}\left(\left[\varphi_{2}\right]\right)\right)$ we have

$$
\operatorname{dim}_{W, F}(\varphi) \geq \operatorname{dim}_{W, S_{F}\left(\varphi_{1}\right)}[\varphi] \geq \operatorname{dim}_{W, S_{F_{1}}\left(\left[\varphi_{2}\right]\right)}\left[\psi_{2}\right]
$$

and the same cases I, II and III for $\psi_{2}$. Suppose without loss of generality that we are in case III, i.e, that $\left[\varphi_{3}\right]+\ldots+\left[\varphi_{r}\right]$ is isotropic in $S_{F_{1}}\left(\left[\varphi_{2}\right]\right)$. If $\left[\varphi_{j}\right]$ is isotropic in $S_{F_{1}}\left(\left[\varphi_{2}\right]\right)$ for all $j \geq 3$, then we are in case I. Now suppose $\left[\varphi_{3}\right]$ anisotropic in $S_{F_{1}}\left(\left[\varphi_{2}\right]\right)$ and denote $F_{2}:=S_{F_{1}}\left(\left[\varphi_{2}\right]\right)$. In $S_{F_{2}}\left(\left[\varphi_{3}\right]\right)$ (which is a special hyperfield) look at

$$
\psi_{3}:=\left[\left[\varphi_{3}\right]+\ldots+\left[\varphi_{r}\right]\right]=2^{n}\langle 1\rangle+\left[\varphi_{4}\right] \ldots+\left[\varphi_{r}\right] \in S_{F_{2}}\left(\left[\varphi_{3}\right]\right) .
$$

For $\psi_{3} \in I^{n}\left(S_{F_{2}}\left(\left[\varphi_{3}\right]\right)\right)$ we have

$$
\operatorname{dim}_{W, F}(\varphi) \geq \operatorname{dim}_{W, S_{F}\left(\varphi_{1}\right)}[\varphi] \geq \operatorname{dim}_{W, S_{F_{1}}\left(\left[\varphi_{2}\right]\right)}\left[\psi_{2}\right] \geq \operatorname{dim}_{W, S_{F_{2}}\left(\left[\varphi_{3}\right]\right)}\left[\psi_{3}\right] .
$$

and the same cases I, II and III for $\psi_{3}$. Repeating this process more $r-3$ times, we get at $\left[\varphi_{r}\right]$ in $S_{F_{r-1}}\left(\left[\varphi_{r-1}\right]\right)$ and

$$
\begin{aligned}
\operatorname{dim}_{W, F}(\varphi) & \geq \operatorname{dim}_{W, S_{F}\left(\varphi_{1}\right)}[\varphi] \geq \operatorname{dim}_{W, S_{F_{1}}\left(\left[\varphi_{2}\right]\right)}\left[\psi_{2}\right] \\
& \geq \operatorname{dim}_{W, S_{F_{2}}\left(\left[\varphi_{3}\right]\right)}\left[\psi_{3}\right] \geq \ldots \geq \operatorname{dim}_{W, S_{F_{r-2}}\left(\left[\varphi_{r-2}\right]\right)}\left[\psi_{r-1}\right] .
\end{aligned}
$$

Now, if $\left[\varphi_{r}\right]$ is isotropic in $S_{F_{r-1}}\left(\left[\varphi_{r-1}\right]\right)$ then $\left[\varphi_{r}\right]$ is hyperbolic in $S_{F_{r-1}}\left(\left[\varphi_{r-1}\right]\right)$, which by case I imply $\operatorname{dim}_{W, S_{F_{r-2}}\left(\left[\varphi_{r-2}\right]\right)}\left[\psi_{r-1}\right] \geq 2^{n}$. If $\left[\varphi_{r}\right]$ is anisotropic in $S_{F_{r-1}}\left(\left[\varphi_{r-1}\right]\right)$ we are in case II and also $\operatorname{dim}_{W, S_{F_{r-2}}\left(\left[\varphi_{r-2}\right]\right)}\left[\psi_{r-1}\right] \geq 2^{n}$.

Now suppose $-1 \in D_{F}(\varphi)$. Then $S_{F}\left(\varphi_{1}\right) \cong\{0,1\}$ (see Theorem 6.2.1], which imply [ $\varphi$ ] is hyperbolic, enabling us to use the very an adapted version argument in Case (I) above: the equivalence class of $\varphi$ on $S_{F}\left(\varphi_{1}\right)$ still is given by

$$
[\varphi]=\left[\varphi_{1}+\ldots+\varphi_{r}\right]=\left[\varphi_{1}\right]+\ldots+\left[\varphi_{r}\right]=2^{n}\langle 1\rangle+\left[\varphi_{2}\right]+\ldots+\left[\varphi_{r}\right] .
$$

Then we have $\left[\varphi_{2}\right]+\ldots+\left[\varphi_{r}\right] \equiv_{S_{F}\left(\varphi_{1}\right)} 2^{n}\langle-1\rangle$, implying that

$$
\operatorname{dim}_{W, F}(\varphi) \geq \operatorname{dim}_{W, F}\left(\varphi_{2}+\ldots+\varphi_{r}\right) \geq \operatorname{dim}_{W, F}\left(\left[\varphi_{2}\right]+\ldots+\left[\varphi_{r}\right]\right)=2^{n}
$$

Now, we turn our attention to graded rings associated to abstract quadratic forms theories (special hyperfields, or equivalently, special groups): we will apply the above established Theorem APH to obtain information on the inductive graded rings (Definition 3.1 in [27]) of a special group $G$ : the graded Witt ring of $G$,

$$
W_{*}(G)=\left(I^{n}(G) / I^{n+1}(G) \xrightarrow{\langle 1,1\rangle \otimes-} I^{n+1}(G) / I^{n+2}(G)\right)_{n \in \mathbb{N}},
$$

and on the graded ring of $k$-theory of $G$,

$$
k_{*}(G)=\left(k_{n}(G) \xrightarrow{\lambda(-1) \otimes-} k_{n+1}(G)\right)_{n \in \mathbb{N}} .
$$

The uses of k-theoretic (and Boolean) methods in abstract theories of quadratic forms has been proved a very successful method, see for instance, these two papers of Dickmann and Miraglia:
[27] where they give an affirmative answer to Marshall's Signature Conjecture, and [29], where they give an affirmative answer to Lam's Conjecture (previously both conjecture have kept open for almost three decades). These two central papers makes us take a deeper look at the theory of special groups (and hence, hyperbolic/pre-special hyperfields) by itself. This is not mere exercise in abstraction: from Marshall's and Lam's Conjecture many questions arise in the abstract and concrete context of quadratic forms.

We will freely permute between a special group $G$ and a special hyperfield $F$ since the associations $G \mapsto F_{G}:=G \dot{\cup}\{0\}$ and $F \mapsto G_{F}:=F \backslash\{0\}$ are part of an equivalence of categories ([24], [17]). The graded Witt ring of a special group is studied in [27] and [28]; [30] is the reference for the k-theory of special groups; in [18] is developed a k-theory for all hyperbolic hyperfields (that includes all pre-special hyperfields).

For the reader's convenience we recall below some relevant definitions.
defn:ksg-aph
Definition 6.3.3 (The Dickmann-Miraglia k-theory [30]). For each special group $G$ (written multiplicatively) we associate a (inductive) graded ring

$$
k_{*} G=\left(k_{0} G, k_{1} G, \ldots, k_{n} G, \ldots\right)
$$

as follows: $k_{0} G:=\mathbb{F}_{2}$ and $k_{1} G:=G$ written additively. With this purpose, we fix the canonical "logarithm" isomorphism $\lambda: G \rightarrow k_{1} G, \lambda(a b)=\lambda(a)+\lambda(b)$. Observe that $\lambda(1)$ is the zero of $k_{1} G$ and $k_{1} G$ has exponent 2, i.e, $\lambda(a)=-\lambda(a)$ for all $a \in G$. In the sequel, we define $k_{*} G$ by the quotient of the $\mathbb{F}_{2}$-graded algebra

$$
\left(\mathbb{F}_{2}, k_{1} G, k_{1} G \otimes_{\mathbb{F}_{2}} k_{1} G, k_{1} G \otimes_{\mathbb{F}_{2}} k_{1} G \otimes_{\mathbb{F}_{2}} k_{1} G, \ldots\right)
$$

by the (graded) ideal generated by $\left\{\lambda(a) \otimes \lambda(a b), a \in D_{G}(1, b)\right\}$. In other words, for each $n \geq 2$,

$$
k_{n} G:=T^{n}\left(k_{1} G\right) / Q^{n}(G),
$$

where

$$
T^{n}\left(k_{1} G\right):=k_{1} G \otimes_{\mathbb{F}_{2}} k_{1} G \otimes_{\mathbb{F}_{2}} \ldots \otimes_{\mathbb{F}_{2}} k_{1} G
$$

and $Q^{n}(G)$ is the subgroup generated by all expressions of type $\lambda\left(a_{1}\right) \otimes \lambda\left(a_{2}\right) \otimes \ldots \otimes \lambda\left(a_{n}\right)$ such that for some $i$ with $1 \leq i<n$, there exist $b \in G$ such that $a_{i} \in D_{G}(1, b)$ and $a_{i}=a_{i+1} b$, which in symbols, means

$$
\begin{aligned}
Q^{n}(G) & :=\left\langle\left\{\lambda\left(a_{1}\right) \otimes \lambda\left(a_{2}\right) \otimes \ldots \otimes \lambda\left(a_{n}\right): \text { exists } 1 \leq i<n \text { and } b \in G\right.\right. \\
& \text { such that } \left.\left.a_{i}=a_{i+1} b \text { and } a_{i} \in D_{G}(1, b)\right\}\right\rangle .
\end{aligned}
$$

Definition 6.3.4 ([27], 30]). Let $G$ be a formally real special group.
$a$ - It holds $[M C(G)]$ (i.e., $G$ satisfies "Marshall's conjecture") if for all $n \geq 1$ and all forms $\varphi$ over $G$,

$$
\text { For all } \sigma \in X_{G} \text {, if } \sigma(\varphi) \equiv 0 \bmod 2^{n} \text { then } \varphi \in I^{n} G \text {. }
$$

$b$ - It holds $[W M C(G)]$ (i.e., $G$ satisfies "Weak Marshall's conjecture") if for all $n \geq 1$, the multiplication by $\langle 1,1\rangle$ is an injection of $I^{n}(G) / I^{n+1}(G)$ into $I^{n+1}(G) / I^{n+2}(G)$.
$c$ - It holds $[S M C(G)]$ (i.e., $G$ satisfies "Strong Marshall's conjecture") if for all $n \geq 1$, the multiplication by $\lambda(-1)$ is an injection of $k_{n}(G)$ into $k_{n+1}(G)$.

It follows from Proposition 4.6.(e) in [27] that $[\mathrm{MC}(G)] \Rightarrow[\mathrm{WMC}(G)]$; in Proposition 4.4 in [29] is established $[\operatorname{SMC}(G)] \Rightarrow[\mathrm{MC}(G)]$, for all reduced special group $G$. Now we apply Theorem APH to obtain the following:

MC-aph

Proposition 6.3.5. Let $G$ be a formally real special group. Then $G$ satisfy Marshall [MC] (i.e., Marshall's signature conjecture holds in $G$ ) iff $G$ satisfy [WMC] (i.e., Weak Marshall's conjecture holds in $G$ ).

Proof. In the theorem 5.3 of [27] is established the equivalence of [MC] and [WMC] for all formally real special groups $G$ such that $2^{k}=\langle 1,1\rangle^{k} \notin I^{k+1}(G)$,for all $k \geq 1$. But, it follows from Theorem APH that all formally real special groups automatically satisfies that property: otherwise $\langle 1,1\rangle^{k}$ will be hyperbolic and thus $-1 \in \operatorname{Sat}(G)=\bigcup_{k \in \mathbb{N}} D_{G}\left(2^{k}\right)$, contradicting that $G$ is a formally real special group.

```
igr1-aph
```

Definition 6.3.6. An inductive graded ring (or Igr for short) is a structure $\mathcal{R}=\left(\left(R_{n}\right)_{n \geq 0},\left(h_{n}\right)_{n \geq 0}, *_{n m}\right)$ where

$$
\begin{aligned}
& i-R_{0} \cong \mathbb{F}_{2} . \\
& i i-R_{n} \text { is a group of exponent } 2 \text { with a distinguished element } \top_{n} \text {. }
\end{aligned}
$$

iii $-h_{n}: R_{n} \rightarrow R_{n+1}$ is a group homomorphism such that $h_{n}\left(T_{n}\right)=T_{n+1}$.
iv - For all $n \geq 0, h_{n}=*_{1 n}\left(\top_{1},-\right)$.
$v$ - The ring

$$
R=\bigoplus_{n \geq 0} R_{n}
$$

is a commutative graded ring.
vi- For $0 \leq s \leq t$ define

$$
h_{s}^{t}=\left\{\begin{array}{l}
I d_{R_{s}} \text { if } s=t \\
h_{t-1} \circ \ldots \circ h_{s+1} \circ h_{s} \text { if } s<t .
\end{array}\right.
$$

Then if $p \geq n$ and $q \geq m$, for all $x \in R_{n}$ and $y \in R_{m}$,

$$
h_{n}^{p}(x) * h_{m}^{q}(y)=h_{n+m}^{p+q}(x * y) .
$$

A morphism between Igr's $\mathcal{R}$ and $\mathcal{R}^{\prime}$ is a morphism of pointed groups and

$$
f=\bigoplus_{n \geq 0} f_{n}: R \rightarrow R^{\prime}
$$

is a morphism of commutative rings with unity (thus $\alpha_{n+1} \circ h_{n}=h_{n+1}^{\prime} \circ \alpha_{n}$ ). The category of inductive graded rings (in first version) and their morphisms will be denoted by IGR.

In [18] are considered some full subcategories of $I G R$ and [21] deals with limits and colimits of $I G R$ and these subcategories. A particularly useful subcategories is $I G R_{h}$, the full subcategory of Igr's $\mathcal{R}$ where for each $a \in R_{1}, \top_{1} *_{1,1} a=a *_{1,1} a \in R_{2}$. Proposition 4.18 and Definition 4.19 therein describes a functor $\Gamma: I G R_{h} \rightarrow p S G$ (the category of pre-special groups, that is equivalent to the category of pre-special hyperfields).

## 174CHAPTER 6. QUADRATIC EXTENSIONS OF SPECIAL GROUPS, HAUPTSATZ AND CONSEQUEN

Now, let $R \in I G R_{h}$. We have a pre-special group $\Gamma(\mathcal{R})=(G(\mathcal{R}),+,-., 0,1)$ by the following: firstly, fix an isomorphism $e_{R}:\left(R_{1},+_{1}, 0_{1}, T_{1}\right) \rightarrow(G(\mathcal{R}), \cdot, 1,-1)$. This isomorphism makes, for example, an element $a *_{11}\left(T_{1}+b\right) \in R_{2}, a, b \in R_{1}$ take the form $\left(e_{R}^{-1}(x)\right) *_{11}\left(e_{R}^{-1}((-1) \cdot y)\right) \in R_{2}$, $x, y \in G(\mathcal{R})$.

Now, let $\Gamma(\mathcal{R}):=G(\mathcal{R})$ and for $a, b, c, d \in R_{1}$ we have $\left\langle e_{R}(a), e_{R}(b)\right\rangle \equiv\left\langle e_{R}(c), e_{R}(d)\right\rangle$ iff $a+b=c+d \in R_{1}$ and $a *_{11} b=c *_{11} d \in R_{2}$

If $\alpha=\left(\alpha_{n}\right)_{n \in \mathbb{N}}: \mathcal{R} \rightarrow \mathcal{R}^{\prime}$ is a IGR-morphism, then $\Gamma(\alpha): G(\mathcal{R}) \rightarrow G\left(\mathcal{R}^{\prime}\right)$ is the unique function (that turns out to be a pSG-morphism) such that $\Gamma(\alpha)=e_{R^{\prime}} \circ \alpha_{1} \circ e_{R}^{-1}$.

For each $G \in p S G$, the Igr's $W_{*}(G)$ and $k_{*}(G)$ belongs to the subcategory $I G R_{h}$ (Lemma 3.2 in [27], [30] and Lemma 9.12 in [28]) is defined a IGR-morphism $s_{G}: k_{*}(G) \rightarrow W_{*}(G)$ such that: $\left(s_{G}\right)_{n}\left(\lambda\left(a_{1}\right) \otimes \cdots \otimes \lambda\left(a_{n}\right)\right)=\left\langle 1,-a_{1}\right\rangle \otimes \cdots \otimes\left\langle 1,-a_{n}\right\rangle$ (Theorem 4.1 in [29]). In general $\left(s_{G}\right)_{n}:$ $k_{n}(G) \rightarrow I^{n}(G) / I^{n+1}(G)$ is a surjective homomorphism of pointed 2-groups and if $n=0,1,2$, then $\left(s_{G}\right)_{n}$ is an isomorphism of pointed 2-groups.

Theorem 4.20 in [18] establishes that the functor $k_{*}: p S G \rightarrow I G R_{q}$ is left adjoint to the functor $\Gamma$ and the natural transformation that is the unity of this adjunction, $\kappa=\left(\kappa_{G}\right)_{G \in p S G}$, is such that for each $G \in p S G, \kappa_{G}: G \rightarrow \Gamma\left(k_{*}(G)\right), g \mapsto \lambda(g)$ is a pSG-morphism that is an isomorphism of the underlying pointed 2 -groups.

In [46] M. Marshall proved that $\omega_{G}: G \rightarrow I(G) / I^{2}(G) g \mapsto\langle 1,-g\rangle+I^{2}(G)$ is an isomorphism of pointed groups such that for each $a, b, c, d \in G$ :

$$
\langle a, b\rangle \equiv_{G}\langle c, d\rangle \Rightarrow\langle 1,-a\rangle \otimes\langle 1,-b\rangle+I^{3}(G)=\langle 1,-c\rangle \otimes\langle 1,-d\rangle+I^{3}(G) .
$$

Thus, for each special group $G$, we have the following commutative diagram of pre-special groups and pSG-morphisms

$$
\left(G \xrightarrow{\omega_{G}} \Gamma\left(W_{*}(G)\right)\right)=\left(G \xrightarrow{\kappa_{G}} \Gamma\left(k_{*}(G)\right) \xrightarrow{\Gamma\left(s_{G}\right)} \Gamma\left(W_{*}(G)\right)\right)
$$

that is, moreover, natural in $G$. Now we are in position to state the:
k-stable-aph
Proposition 6.3.7. Let $G$ be a special group. Then

$$
\left.i-\Gamma\left(s_{G}\right): \Gamma\left(k_{*}(G)\right) \rightarrow \Gamma\left(W_{*}(G)\right)\right) \text { is a } p S G \text {-isomorphism. }
$$

ii $-\omega_{G}: G \rightarrow \Gamma\left(W_{*}(G)\right)$ is a $p S G$-isomorphism.
iii - $\kappa_{G}: G \rightarrow \Gamma\left(k_{*}(G)\right)$ is a $p S G$-isomorphism.
In particular, $\Gamma\left(k_{*}(G)\right)$ and $\Gamma\left(W_{*}(G)\right)$ are special groups.
Proof. This is essentially contained in the proof of Lemma 3.5 in [30], but for convince the reader we provide some details:

First observe that, from axiom [SG4] is enough to show that for $x, y \in G$ are equivalent:
(1) $x \in D_{G}(1, y)$;
(2) $\lambda(x) \in D_{\Gamma\left(k_{*}(G)\right)}(\lambda(1), \lambda(y))$;
(3) $\langle 1,-x\rangle+I^{2}(G) \in D_{\Gamma\left(W_{*}(G)\right)}\left(\langle 1,-1\rangle+I^{2}(G),\langle 1,-y\rangle+I^{2}(G)\right)$
$(1) \Rightarrow(2)$ : is clear from the definition of $k_{*}(G)$, just note that condition (2) is equivalent to $\lambda(x) \lambda(x y)=0 \in k_{2}(G)$
$(2) \Rightarrow(3)$ : this follows directly from $s_{G}: k_{*}(G) \rightarrow W_{*}(G)$ be a $I G R_{h}$-morphism
$(3) \Rightarrow(1)$ : note that condition (3) is equivalent to $\left.\left.\langle 1,-x\rangle \otimes\langle 1,-x y\rangle+I^{3}(G)\right)=0+I^{3}(G)\right)$. This means that $\langle 1,-x\rangle \otimes\langle 1,-x y\rangle \in I^{3}(G)$. Since $\operatorname{dim}(\langle 1,-x\rangle \otimes\langle 1,-x y\rangle)=4<8=2^{3}$, then by Theorem APH, $\langle 1,-x\rangle \otimes\langle 1,-x y\rangle$ is an isotropic Pfister form. Thus it is an hyperbolic form, and then, by Proposition 2.2.(k) in [28], $x \in D_{G}(1, y)$.

The notion of k-stable hyperbolic hyperfield $F$, i.e. those such that canonical morphism $\kappa_{F}$ : $F \rightarrow \Gamma\left(k_{*}(F)\right) \dot{\cup}\{0\}$ is an isomorphism of hyperfields, it is fundamental in [14]. Thus the previous result establishes the:
k-stable-co

Corollary 6.3.8. Every special hyperfield $F$ is $k$-stable.
The following result shows that the k -theory construction provides a very good encoding of the neither complete neither cocomplete - category of special groups into the complete and cocomplete category of inductive graded rings.
epiK-aph

Proposition 6.3.9. The functor $k_{*}: S G \rightarrow I G R$ is full and faithful.
Proof. We have to show that for each special groups $G_{0}$ and $G_{1}$ and any $\beta: k_{*} G_{0} \rightarrow k_{*} G_{1}$ be an inductive graded ring morphism between the associated inductive graded rings of $k$-theory, then there exist unique SG-morphism $f: G_{0} \rightarrow G_{1}$ such that $\beta=k_{*}(f)$.

Proposition 3.6 in [30] establishes (from Lemma 3.5) that: for each special groups $G_{0}$ and $G_{1}$ and any $\beta: k_{*}\left(G_{0}\right) \rightarrow k_{*}\left(G_{1}\right)$ be a graded ring morphism between the induced $k$-theory graded rings, such that $\beta_{0}=i d_{\mathbb{F}_{2}}$ and $G_{1}$ is a $\mathrm{AP}(3)$ special group, then there exist a qSG-morphism $f: G \rightarrow H$ such that $\beta=k_{*}(f)$. Moreover, this $f$ is uniquely determined since $\beta_{1} \circ \lambda_{G_{0}}=\lambda_{G_{1}} \circ f$ and $\lambda_{G_{i}}: G_{i} \rightarrow k_{1}\left(G_{i}\right)$ is an isomorphism of groups of exponent 2 that preserves the distinguished elements $\left(-1_{G_{i}} \mapsto \lambda_{G_{i}}\left(-1_{G_{i}}\right)\right.$.

Thus the result follows since any special group $G_{1}$ satisfies is $\mathrm{AP}(3)$ (by Theorem APH), and since $\beta$ is a IGR-morphism then automatically $\beta_{0}=i d_{\mathbb{F}_{2}}$ and $\beta_{1}=k_{1}(f)$ implies that $f\left(-1_{G_{0}}\right)=$ $-1_{G_{1}}$, thus $f$ the qSG-morphism $f$ is a SG-morphism.

Remark 6.3.10. The previous result can be derived, alternatively from Corollary 6.3 .8 and Theorem 4.20 in [18]: from an well known result on adjoint functors, a left adjoint is a full and faithful functor iff the unity of the adjunction is an isomorphism. Thus $k_{*}: S G \rightarrow I G R_{h}$ is a full and faithful functor and, since $I G R_{h} \subseteq I G R$ is a full subcategory, then $k_{*}: S G \rightarrow I G R$ is full and faithful.

## Chapter 7

## The Galois group of a Special Group

The Igr's functors $W_{*}, k_{*}$ were extended by M. Dickmann and F.Miraglia from the category of fields of characteristic $\neq 2$ to the category of special groups (equivalently, the category of special hyperfields). Another relevant Igr functor, the graded cohomology ring, $H^{*}\left(\operatorname{Gal}\left(F^{s} \mid F\right),\{ \pm 1\}\right)$ remains defined only on the field setting. This chapter constitutes an attempt to provide an Igr functor associated to a (Galois) cohomology theory for special groups, based on the work of J. Minac and M. Spira [53]. We will define - by "generator and relations", Gal(G), the Galois Group of an $S G G$ (Definition 7.2.10), and provide some properties of this construction, as the encoding of the orderings on $G$. However, since deeper results will depend of a description of $G a l(G)$ "from below", and it still unavailable a complete theory of algebraic extension of (super)hyperfields, we will not pursue a more complete development of this cohomology theory in this thesis, reserving it for a future research. The main results here, established for the "standard pre-special groups", are Theorems $7.3 .12,7.3 .13$ and 7.3 .15 , that recover for the abstract context the characterization of orderings in terms of the involutions in the Galois group of a field. These results provide a clue that this profinite group $G a l(G)$, defined by generators and relations, is not -at this moment- a legitimate Galois group (since a characterization of it from quadratic subextensions is still missing), but it works in some aspects as a Galois group of a field in the sense that it can encode faithfully some relevant information on the structure of $G$.

We will work in the category of pro-2-groups, and take, as usual, the conventions: "subgroup" means "closed subgroup"; "subgroup generated by a subset" means "the closure of the abstract group generated by the subset"; "morphism" means "continuous homomorphism".

### 7.1 The motivation: W-groups

The context that we will keep in mind is essentially that of the results developed in sections 1 and 2 of [53]. In this Section we will reproduce (and expand the details of) part of these results.

Consider a field $F$. In [53], J. Minac and M. Spira define a special Galois extension of the base field $F$, and determine its structure and its Galois group through the behavior of quaternions algebras over $F$. As they developed in [53], this extension contain essentially all the information need to understand the behavior of quadratic forms over $F$.

Recall that the quadratic closure of $F$, denoted by $F_{q}$, is the smallest extension of $F$ which is closed under taking of square roots (or more explicit, the compositum of all 2-towers over $F$ inside a fixed algebraic closure of $F)$. The group $\operatorname{Gal}_{F}\left(F_{q}\right)$ will be denoted by $G_{F}^{q}$.

Let $\left\{a_{i}: i \in I\right\}$ be a basis of $\dot{F} / \dot{F}^{2}$ (as Minac and Spira did in their paper, we will assume that $1,2, . ., n \in I$ with $1<2<\ldots<n$ in order to easy our presentation). We define $F^{(2)}=$
$F(\sqrt{a}: a \in \dot{F})$ (note that $F^{(2)}=F\left(\sqrt{a_{i}}: i \in I\right), \mathcal{E}=\left\{y \in F^{(2)}: F^{(2)}(\sqrt{y}) \mid F\right.$ is Galois $\}$ and $F^{(3)}=F^{(2)}(\sqrt{y}: y \in \mathcal{E})$ if $F$; if $F$ is quadratically closed we set $F^{(2)}=F^{(3)}=F$.

Minac and Spira built a strong connection between $F^{(3)}$ and the Witt ring of $F$. They named $F^{(3)}$ as Witt closure of $F$. The group $G_{F}:=\operatorname{Gal}_{F}\left(F^{3} \mid F\right)$ is called the $\mathbf{W}$-group of $F$. Our goal here is to describe a way to factor $G_{F}$ as $G_{F} \cong \mathcal{W}(I) / \mathcal{V}(I)$, with $\mathcal{W}(I)$ and $\mathcal{V}(I)$ interesting profinite groups. This procedure will reveal how to generalize $G_{F}$ in the context of abstract theories of quadratic forms, and in particular, describe what would be a Galois group associated to a special group. The first step is to describe $\mathcal{W}(I)$.

For an arbitrary group $G$, define $\hat{G}=G^{4}\left[G^{2}, G\right]$, i.e, the (closed) subgroup generated by fourth powers and by commutators of the form $\left[g^{2}, h\right]$ for $g, h \in G$. Let $t^{4}\left[g^{2}, h\right] \in \hat{G}$. Then, for each $z \in G$ :

$$
z^{-1}\left(t^{4}\left[g^{2}, h\right]\right) z=z^{-1} t^{4}\left[g^{2}, h\right] z=\left(z^{-1} t^{4} z\right)\left(z^{-1}\left[g^{2}, h\right] z\right)
$$

with $z^{-1} t^{4} z=\left(z^{-1} t z\right)^{4} \in G^{4}$ and $z^{-1}\left[g^{2}, h\right] z \in\left[G^{2}, G\right]$, because

$$
\begin{aligned}
z^{-1}\left[g^{2}, h\right] z & =z^{-1} g^{-2} h^{-1} g^{2} h z \\
& =\left(z^{-1} g^{-2} z\right)\left(z^{-1} h^{-1} z\right)\left(z^{-1} g^{2} z\right)\left(z^{-1} h z\right) \\
& =\left(z^{-1} g z\right)^{-2}\left(z^{-1} h z\right)^{-1}\left(z^{-1} g z\right)^{2}\left(z^{-1} h z\right) \\
& =\left[\left(z^{-1} g z\right)^{2},\left(z^{-1} h z\right)\right] .
\end{aligned}
$$

Hence $\hat{G}$ is a normal subgroup of $G$, and we define $\bar{G}=G / \hat{G}$. Let $\mathcal{C}$ denote the class of profinite 2 -groups $G$ such that $\hat{G}=\{1\}$.

The main example (and the motivation to consider this full subcategory of pro-2-groups) is the following fact: If $\operatorname{char}(F) \neq 2$ then $G=G a l\left(F^{(3)} \mid F\right)$ satisfies this condition $\hat{G}=\{1\}$, since, by Proposition 2.1 in [53], $G \cong G_{q} / G_{q}{ }^{4} \cdot\left[G_{q}{ }^{2}, G_{q}\right]$, where $G_{q}=G a l\left(F_{q} \mid F\right)$ and $F_{q}$ is a quadratic closure of $F$.

A pro- $\mathcal{C}$-group will be called just $\mathcal{C}$-group, and $\mathcal{C}$-group on $I$ if it has a minimal set of generators of cardinality $|I|$.

Let $I$ be a well-ordered set. The next step is to describe a canonical way to represent the elements of $\bar{S}$ where $S$ is the free pro-2-group on a nonempty set $I$. Let

$$
\mathcal{W}(I):=\prod_{i \in I} \mathbb{Z}_{2} \times \prod_{\substack{i, j \in I \\ i<j}} \mathbb{Z}_{2} \times \prod_{i \in I} \mathbb{Z} / 2 \mathbb{Z}
$$

Here we are considering $\mathbb{Z}_{2}$ multiplicatively, i.e, $\mathbb{Z}_{2} \cong\{1,-1\}$. A typical element of $\mathcal{W}_{I}$ will be written as $\left(t_{i}^{\alpha_{i}}\right)\left(t_{i j}^{\beta_{i j}}\right)\left(x_{i}^{\gamma_{i}}\right)$, where $\alpha_{i}, \beta_{i j}, \gamma_{i} \in\{0,1\}$ and $t_{i}=t_{i j}=x_{i}=-1$ for all $i, j \in I$.

Let $g, h \in \mathcal{W}(I)$

$$
\begin{aligned}
g & =\left(t_{i}^{\alpha_{i}}\right)\left(t_{i j}^{\beta_{i j}}\right)\left(x_{i}^{\gamma_{i}}\right) \\
h & =\left(t_{i}^{\alpha_{i}^{\prime}}\right)\left(t_{i j}^{\beta_{i j}^{\prime}}\right)\left(x_{i}^{\gamma_{i}^{\prime}}\right)
\end{aligned}
$$

We define

$$
g h=\left(t_{i}^{\alpha_{i}+\alpha_{i}^{\prime}+\gamma_{i} \gamma_{i}^{\prime}}\right)\left(t_{i j}^{\beta_{i j}+\beta_{i j}^{\prime}+\gamma_{i}^{\prime} \gamma_{j}}\right)\left(x_{i}^{\gamma_{i}+\gamma_{i}^{\prime}}\right),
$$

where the exponents are taken modulo 2 .
Lemma 7.1.1. With the notation described above, $\mathcal{W}(I)$ is a group.

## Proof. Let

$$
\begin{aligned}
g & =\left(t_{i}^{\alpha_{i}}\right)\left(t_{i j}^{\beta_{i j}}\right)\left(x_{i}^{\gamma_{i}}\right) \\
h & =\left(t_{i}^{\alpha_{i}^{\prime}}\right)\left(t_{i j}^{\beta_{i j}^{\prime}}\right)\left(x_{i}^{\gamma_{i}^{\prime}}\right) \\
k & =\left(t_{i}^{\alpha_{i}^{\prime \prime}}\right)\left(t_{i j}^{\beta_{i j}^{\prime \prime}}\right)\left(x_{i}^{\gamma_{i}^{\prime \prime}}\right) \\
1 & =\left(t_{i}^{0}\right)\left(t_{i j}^{0}\right)\left(x_{i}^{0}\right)
\end{aligned}
$$

First of all,

$$
\begin{aligned}
& g \cdot 1=\left(t_{i}^{\alpha_{i}+0+\gamma_{i} \cdot 0}\right)\left(t_{i j}^{\beta_{i j}+0+0 \cdot \gamma_{j}}\right)\left(x_{i}^{\gamma_{i}+0}\right)=g \\
& 1 \cdot g=\left(t_{i}^{0+\alpha_{i}+0 \cdot \gamma_{i}}\right)\left(t_{i j}^{0+\beta_{i j}+\gamma_{i} \cdot 0}\right)\left(x_{i}^{0+\gamma_{i}}\right)=g
\end{aligned}
$$

hence 1 is in fact the neutral element of $\mathcal{W}(I)$. In order to find $g^{-1}$, we need to solve a system of modulo 2 congruences.

$$
g h:=1 \Rightarrow\left\{\begin{array} { l } 
{ \alpha _ { i } + \alpha _ { i } ^ { \prime } + \gamma _ { i } \gamma _ { i } ^ { \prime } \equiv _ { 2 } 0 } \\
{ \beta _ { i j } + \beta _ { i j } ^ { \prime } + \gamma _ { i } ^ { \prime } \gamma _ { j } \equiv _ { 2 } 0 } \\
{ \gamma _ { i } + \gamma _ { i } ^ { \prime } \equiv _ { 2 } 0 }
\end{array} \Rightarrow \left\{\begin{array}{l}
\gamma_{i}^{\prime} \equiv_{2} \gamma_{i} \\
\beta_{i j}^{\prime} \equiv_{2} \beta_{i j}+\gamma_{i} \gamma_{j} \\
\alpha_{i}^{\prime} \equiv_{2} \alpha_{i}+\gamma_{i}
\end{array}\right.\right.
$$

Then, taking

$$
g^{-1}=\left(t_{i}^{\alpha_{i}+\gamma_{i}}\right)\left(t_{i j}^{\beta_{i j}+\gamma_{i} \gamma_{j}}\right)\left(x_{i}^{\gamma_{i}}\right)
$$

we obtain $g \cdot g^{-1}=g^{-1} \cdot g=1$.
Finally, for associativity, we have

$$
\begin{aligned}
(g \cdot h) \cdot k & =\left[\left(t_{i}^{\alpha_{i}+\alpha_{i}^{\prime}+\gamma_{i} \gamma_{i}^{\prime}}\right)\left(t_{i j}^{\beta_{i j}+\beta_{i j}^{\prime}+\gamma_{i}^{\prime} \gamma_{j}}\right)\left(x_{i}^{\gamma_{i}+\gamma_{i}^{\prime}}\right)\right] \cdot k \\
& =\left(t_{i}^{\alpha_{i}+\alpha_{i}^{\prime}+\alpha_{i}^{\prime \prime}+\gamma_{i} \gamma_{i}^{\prime}+\left(\gamma_{i}+\gamma_{i}^{\prime}\right) \gamma_{i}^{\prime \prime}}\right)\left(t_{i j}^{\beta_{i j}+\beta_{i j}^{\prime}+\beta_{i j}^{\prime \prime}+\gamma_{i}^{\prime} \gamma_{j}+\gamma_{i}^{\prime \prime}\left(\gamma_{j}+\gamma_{j}^{\prime}\right)}\right)\left(x_{i}^{\gamma_{i}+\gamma_{i}^{\prime}+\gamma_{i}^{\prime \prime}}\right) \\
& =\left(t_{i}^{\alpha_{i}+\alpha_{i}^{\prime}+\alpha_{i}^{\prime \prime}+\gamma_{i} \gamma_{i}^{\prime}+\gamma_{i} \gamma_{i}^{\prime \prime}+\gamma_{i}^{\prime} \gamma_{i}^{\prime \prime}}\right)\left(t_{i j}^{\beta_{i j}+\beta_{i j}^{\prime}+\beta_{i j}^{\prime \prime}+\gamma_{i}^{\prime} \gamma_{j}+\gamma_{i}^{\prime \prime} \gamma_{j}+\gamma_{i}^{\prime \prime} \gamma_{j}^{\prime}}\right)\left(x_{i}^{\gamma_{i}+\gamma_{i}^{\prime}+\gamma_{i}^{\prime \prime}}\right)
\end{aligned}
$$

and

$$
\begin{aligned}
g \cdot(h \cdot k) & =g \cdot\left[\left(t_{i}^{\alpha_{i}^{\prime}+\alpha_{i}^{\prime \prime}+\gamma_{i}^{\prime} \gamma_{i}^{\prime \prime}}\right)\left(t_{i j}^{\beta_{i j}^{\prime}+\beta_{i j}^{\prime \prime}+\gamma_{i}^{\prime \prime} \gamma_{j}^{\prime}}\right)\left(x_{i}^{\gamma_{i}^{\prime}+\gamma_{i}^{\prime \prime}}\right)\right] \\
& =\left(t_{i}^{\alpha_{i}+\alpha_{i}^{\prime}+\alpha_{i}^{\prime \prime}+\gamma_{i}^{\prime} \gamma_{i}^{\prime \prime}+\gamma_{i}\left(\gamma_{i}^{\prime}+\gamma_{i}^{\prime \prime}\right)}\right)\left(t_{i j}^{\beta_{i j}+\beta_{i j}^{\prime}+\beta_{i j}^{\prime \prime}+\gamma_{i}^{\prime \prime} \gamma_{j}^{\prime}+\left(\gamma_{i}^{\prime}+\gamma_{i}^{\prime \prime}\right) \gamma_{j}}\right)\left(x_{i}^{\gamma_{i}+\gamma_{i}^{\prime}+\gamma_{i}^{\prime \prime}}\right) \\
& =\left(t_{i}^{\alpha_{i}+\alpha_{i}^{\prime}+\alpha_{i}^{\prime \prime}+\gamma_{i}^{\prime} \gamma_{i}^{\prime \prime}+\gamma_{i} \gamma_{i}^{\prime}+\gamma_{i} \gamma_{i}^{\prime \prime}}\right)\left(t_{i j}^{\beta_{i j}+\beta_{i j}^{\prime}+\beta_{i j}^{\prime \prime}+\gamma_{i}^{\prime \prime} \gamma_{j}^{\prime}+\gamma_{i}^{\prime} \gamma_{j}+\gamma_{i}^{\prime \prime} \gamma_{j}}\right)\left(x_{i}^{\gamma_{i}+\gamma_{i}^{\prime}+\gamma_{i}^{\prime \prime}}\right)
\end{aligned}
$$

Thus $(g \cdot h) \cdot k=g \cdot(h \cdot k)$ completing the proof.
Remark 7.1.2. 1. Note that, if $|I|=n$, then

$$
|\mathcal{W}(I)|=\left|\prod_{i \in I} \mathbb{Z}_{2} \times \prod_{\substack{i, j \in I \\ i \leq j}} \mathbb{Z}_{2} \times \prod_{i \in I} \mathbb{Z}_{2}\right|=2^{\left(n^{2}+3 n\right) / 2}
$$

2. The example above is the free $\mathcal{C}$ in $n$-generators. In fact: for each $n \in \mathbb{N}$, let $F(n)$ be the free
group in $n$-generators, then $\mathcal{W}(n) \cong F(n) / \hat{F}(n)$.
3. It follows that the category of finite (discrete) groups $G$ with $\hat{G}=\{1\}$ is a category of $\mathbb{Z}_{2}$ modules that is closed by homomorphic images, subgroups and finite products. In particular, $\mathbb{Z}_{2}=\mathbb{Z} / 2 \mathbb{Z}, \mathbb{Z}_{4}=\mathbb{Z} / 4 \mathbb{Z}$ and $\mathbb{D}_{4} \cong \mathbb{Z}_{2} \ltimes \mathbb{Z}_{4}$ (the 8-element dihedral group), are finite $\mathcal{C}$-groups.
Lets denote, for $k, l \in I$,

$$
\begin{aligned}
t_{k} & :=\left(t_{i}^{\delta_{i k}}\right)\left(1_{i j}\right)\left(1_{i}\right) \\
t_{k l} & :=\left(1_{i}\right)\left(t_{i j}^{\delta_{i j}}(k l)\right)\left(1_{i}\right) \\
x_{k} & :=\left(1_{i}\right)\left(1_{i j}\right)\left(x_{i}^{\delta_{i k}}\right)
\end{aligned}
$$

where for all $i, j, k, l \in I, \delta_{i k}=1$ if $i=k$ and $\delta_{i k}=0$ otherwise; and $\delta_{(i j)(k l)}=1$ if $i=k$ and $j=l$, and $\delta_{(i j)(k l)}=0$ otherwise. After some straightforward calculations we obtain the following results.
Lemma 7.1.3. Consider $t_{k}, t_{k l}, x_{k}$ as above. Then for all $g, h \in \mathcal{W}(I)$, with $g=\left(t_{i}^{\alpha_{i}}\right)\left(t_{i j}^{\beta_{i j}}\right)\left(x_{i}^{\gamma_{i}}\right)$, $h=\left(t_{i}^{\alpha_{i}^{\prime}}\right)\left(t_{i j}^{\beta_{i j}^{\prime}}\right)\left(x_{i}^{\gamma_{i}^{\prime}}\right)$ and $z=\left(t_{i}^{\alpha_{i}^{\prime \prime}}\right)\left(t_{i j}^{\beta_{i j}^{\prime \prime}}\right)\left(x_{i}^{\gamma_{i}^{\prime \prime}}\right)$, we have the following:

$$
i-t_{k} \cdot t_{k}=1
$$

$i i-x_{k} \cdot x_{k}=t_{k}$.
iii- If $k<l$ then $\left[x_{k}, x_{l}\right]=t_{k l}$.
$i v-g^{-1}=\left(t_{i}^{\alpha_{i}+\gamma_{i}}\right)\left(t_{i j}^{\beta_{i j}} \gamma_{i} \gamma_{j}\right)\left(x_{i}^{\gamma_{i}}\right)$.
$v-g^{2}=\left(t_{i}^{\gamma_{i}}\right)\left(t_{i j}^{\gamma_{i} \gamma_{j}}\right)\left(1_{i}\right)$.
$v i-h^{g}=g h g^{-1}=\left(t_{i}^{\alpha_{i}^{\prime}}\right)\left(t_{i j}^{\beta_{i j}^{\prime}+\gamma_{i} \gamma_{j}^{\prime}+\gamma_{i}^{\prime} \gamma_{j}}\right)\left(x_{i}^{\gamma_{i}^{\prime}}\right)$.
$v i i-[g, h]=\left(1_{i}\right)\left(t_{i j}^{\gamma_{i} \gamma_{j}^{\prime}+\gamma_{i}^{\prime} \gamma_{j}}\right)\left(1_{j}\right)$.
viii - $g^{4}=\left[g^{2}, h\right]=1$.
$i x-[[z, w], h]=1$.
Remark 7.1.4. In [53], they simply denote

$$
g=\left(x_{i}^{2 \alpha_{i}}\right)\left(\left[x_{i}, x_{j}\right]^{\beta_{i j}}\right)\left(x_{i}^{\gamma_{i}}\right) .
$$

But we are not going to use this simplification here.
Now, for each $i, j \in I$, consider the following three sets:

$$
\begin{aligned}
M_{i} & :=\left\{g=\left(t_{i}^{\alpha_{i}}\right)\left(t_{i j}^{\beta_{i j}}\right)\left(x_{i}^{\gamma_{i}}\right) \in \mathcal{W}(I): \gamma_{i}=0\right\}, \\
S_{i} & :=\left\{g=\left(t_{i}^{\alpha_{i}}\right)\left(t_{i j}^{\beta_{i j}}\right)\left(x_{i}^{\gamma_{i}}\right) \in \mathcal{W}(I): \alpha_{i}=\gamma_{i}=0\right\}, \\
D_{i j} & :=\left\{g=\left(t_{i}^{\alpha_{i}}\right)\left(t_{i j}^{\beta_{i j}}\right)\left(x_{i}^{\gamma_{i}}\right) \in \mathcal{W}(I): \beta_{i j}=\gamma_{i}=\gamma_{j}=0\right\} .
\end{aligned}
$$

Now, consider the following families

$$
\begin{aligned}
M(I) & :=\left\{M_{i}: i \in I\right\}, S(I):=\left\{S_{i}: i \in I\right\}, D(I):=\left\{D_{i j}: i, j \in I, i \leq j\right\}, \\
V & :=M(I) \cup S(I) \cup D(I) .
\end{aligned}
$$

Proposition 7.1.5. Let $i, j \in I$.
a- $M_{i}$ is a maximal normal subgroup of $\mathcal{W}(I)$ such that $\mathcal{W}(I) / M_{i} \cong \mathbb{Z}_{2}$.
$b-S_{i}$ is a normal subgroup of $\mathcal{W}(I)$ such that $S_{i} \subseteq M_{i}$ and $\mathcal{W}(I) / S_{i} \cong \mathbb{Z}_{4}$.
$c-D_{i j}$ is a normal subgroup of $\mathcal{W}(I)$ such that $D_{i j} \subseteq M_{i} \cap M_{j}$ and $\mathcal{W}(I) / D_{i j} \cong \mathbb{D}_{4}$.
$d-\cap V=\{1\}$.
Proof. We establish the following notation: let $g \in \mathcal{W}(I)$. Then $g=\left(t_{i}^{\alpha_{i}}\right)\left(t_{i j}^{\beta_{i j}}\right)\left(x_{i}^{\gamma_{i}}\right)$ for suitable $\alpha_{i}, \beta_{i j}, \gamma_{i} \in\{0,1\}$. We denote

$$
\alpha_{i}(g):=\alpha_{i}, \beta_{i j}(g):=\beta_{i j} \text { and } \gamma_{i}(g):=\gamma_{i} .
$$

In this sense, if $g, h \in \mathcal{W}(I)$ then

$$
g h=\left(t_{i}^{\alpha_{i}(g)+\alpha_{i}(h)+\gamma_{i}(g) \gamma_{i}(h)}\right)\left(t_{i j}^{\beta_{i j}(g)+\beta_{i j}(h)+\gamma_{i}(h) \gamma_{j}(g)}\right)\left(x_{i}^{\gamma_{i}(g)+\gamma_{i}(h)}\right) .
$$

Moreover, using the formulas in Lemma 7.1.3, we obtain that for all $i, j \in I, M_{i}, S_{i}$ and $D_{i j}$ are proper normal subgroups of $\mathcal{W}(I)$.
a - Let $\tau, \theta \in \mathcal{W}(I) \backslash M_{i}$. Then $\gamma_{i}(\tau)=\gamma_{i}(\theta)=1$ and

$$
\gamma_{i}\left(\theta^{-1} \tau\right)=\gamma_{i}(\theta)+\gamma_{i}(\tau)=0 .
$$

Therefore $\theta^{-1} \tau \in M_{i}$ which imply

$$
\mathcal{W}(I) / M_{i}=\{\overline{1}, \bar{\tau}\} \cong \mathbb{Z}_{2}
$$

b - Note that $S_{i} \subseteq M_{i}$. Now, let $\tau, \theta \in M_{i} \backslash S_{i}$. Then $\alpha_{i}(\tau)=\alpha_{i}(\theta)=1$ and

$$
\alpha_{i}\left(\theta^{-1} \tau\right)=\left[\alpha_{i}(\theta)+\gamma_{i}(\theta)\right]+\alpha_{i}(\tau)+\gamma_{i}(\theta) \gamma_{i}(\tau)=0
$$

Hence $\theta^{-1} \tau \in S_{i}$, and $M_{i} / S_{i} \cong \mathbb{Z}_{2}$. So we have an exact sequence

$$
1 \rightarrow M_{i} / S_{i} \xrightarrow{\iota} \mathcal{W}(I) / S_{i} \xrightarrow{\pi} \mathcal{W}(I) / M_{i} \rightarrow 1,
$$

where $\iota$ and $\pi$ are respectively the canonical inclusion and canonical projection. Moreover

$$
\mathcal{W}(I) / S_{i} \cong M_{i} / S_{i} \rtimes \mathcal{W}(I) / M_{i} \text { or } \mathcal{W}(I) / S_{i} \cong \mathcal{W}(I) / M_{i} \ltimes M_{i} / S_{i} .
$$

In both cases,

$$
\left|\mathcal{W}(I) / S_{i}\right|=\left|M_{i} / S_{i} \times \mathcal{W}(I) / M_{i}\right|=\left|\mathbb{Z}_{2} \times \mathbb{Z}_{2}\right|=4
$$

Now let $\sigma \in \mathcal{M} \backslash M_{i}$. We have $\sigma^{4}=1 \in S_{i}$ with

$$
\gamma_{i}\left(\sigma^{3}\right):=\gamma_{i}(\sigma)=1 .
$$

Then $\bar{\sigma}^{3} \neq \overline{1}$ in $\mathcal{W}(I) / S_{i}$, which proves that $\mathcal{W}(I) / S_{i}$ has an element of order 4 . Then

$$
\mathcal{W}(I) / S_{i} \cong \mathbb{Z}_{4}
$$

c - Remember that

$$
\mathbb{D}_{4}:=\left\langle r, s: r^{4}=s^{2}=(s r)^{2}=1\right\rangle .
$$

Using the same argument of item (b), we get $\left(M_{i} \cap M_{j}\right) / D_{i j} \cong \mathbb{Z}_{2}$ and $\left|\mathcal{W}(I) /\left(M_{i} \cap M_{j}\right)\right|=4$ with

$$
\left|\mathcal{W}(I) / D_{i j}\right|=\left|\left(M_{i} \cap M_{j}\right) / D_{i j} \times \mathcal{W}(I) /\left(M_{i} \cap M_{j}\right)\right|=8 .
$$

More specifically, if we get $\tau_{1} \in\left(M_{i} \cap M_{j}\right) \backslash D_{i j}, \tau_{2}, \theta_{2} \in M_{i} \backslash M_{j}, \tau_{3}, \theta_{3} \in M_{j} \backslash M_{i}$ and $\tau_{4}, \theta_{4} \in \mathcal{W}(I) \backslash\left(M_{i} \cup M_{j}\right)$, with

$$
\begin{aligned}
& \beta_{i j}\left(\tau_{2}\right)=\beta_{i j}\left(\tau_{3}\right)=\beta_{i j}\left(\tau_{4}\right)=1 \\
& \beta_{i j}\left(\theta_{2}\right)=\theta_{i j}\left(\theta_{3}\right)=\theta_{i j}\left(\tau_{4}\right)=0
\end{aligned}
$$

then the following equations hold in $\mathcal{W}(I) / D_{i j}$

$$
\begin{aligned}
& {\overline{\tau_{1}}}^{2}=\overline{1} \\
& \overline{\tau_{4}^{-1} \theta_{4}}=\overline{\tau_{3}^{-1} \theta_{3}}=\overline{\tau_{2}^{-1} \theta_{2}}=\overline{\tau_{1}}
\end{aligned}
$$

Then

$$
\mathcal{W}(I) / D_{i j}=\left\{\overline{1}, \overline{\tau_{1}}, \overline{\tau_{2}}, \overline{\tau_{3}}, \overline{\tau_{4}}, \overline{\tau_{1} \tau_{2}}, \overline{\tau_{1} \tau_{3}}, \overline{\tau_{1} \tau_{4}}\right\},
$$

with the following table of multiplication:

| $\cdot$ | $\overline{\tau_{1}}$ | $\overline{\tau_{2}}$ | $\overline{\tau_{3}}$ | $\overline{\tau_{4}}$ | $\overline{\tau_{1} \tau_{2}}$ | $\overline{\tau_{1} \tau_{3}}$ | $\overline{\tau_{1} \tau_{4}}$ |
| :--- | :--- | :--- | :--- | :--- | :--- | :--- | :--- |
| $\overline{\tau_{1}}$ | $\overline{1}$ | $\overline{\tau_{1} \tau_{2}}$ | $\overline{\tau_{1} \tau_{3}}$ | $\overline{\tau_{1} \tau_{4}}$ | $\overline{\tau_{2}}$ | $\overline{\tau_{3}}$ | $\overline{\tau_{4}}$ |
| $\overline{\tau_{2}}$ | $\overline{\tau_{1} \tau_{2}}$ | $\overline{1}$ | $\overline{\tau_{1} \tau_{4}}$ | $\overline{\tau_{3}}$ | $\overline{\tau_{1}}$ | $\overline{\tau_{4}}$ | $\overline{\tau_{1} \tau_{3}}$ |
| $\overline{\tau_{3}}$ | $\overline{\tau_{1} \tau_{3}}$ | $\overline{\tau_{4}}$ | $\overline{1}$ | $\overline{\tau_{1} \tau_{2}}$ | $\overline{\tau_{1} \tau_{4}}$ | $\overline{\tau_{1}}$ | $\overline{\tau_{2}}$ |
| $\overline{\tau_{4}}$ | $\overline{\tau_{1} \tau_{4}}$ | $\overline{\tau_{1} \tau_{3}}$ | $\overline{\tau_{2}}$ | $\overline{\tau_{1}}$ | $\overline{\tau_{3}}$ | $\overline{\tau_{1} \tau_{2}}$ | $\overline{1}$ |
| $\overline{\tau_{1} \tau_{2}}$ | $\overline{\tau_{2}}$ | $\overline{\tau_{1}}$ | $\overline{\tau_{4}}$ | $\overline{\tau_{1} \tau_{3}}$ | $\overline{1}$ | $\overline{\tau_{1} \tau_{4}}$ | $\overline{\tau_{3}}$ |
| $\overline{\tau_{1} \tau_{3}}$ | $\overline{\tau_{3}}$ | $\overline{\tau_{1} \tau_{4}}$ | $\overline{\tau_{1}}$ | $\overline{\tau_{2}}$ | $\overline{\tau_{4}}$ | $\overline{1}$ | $\overline{\tau_{1} \tau_{2}}$ |
| $\overline{\tau_{1} \tau_{4}}$ | $\overline{\tau_{4}}$ | $\overline{\tau_{3}}$ | $\overline{\tau_{1} \tau_{2}}$ | $\overline{1}$ | $\overline{\tau_{1} \tau_{3}}$ | $\overline{\tau_{2}}$ | $\overline{\tau_{1}}$ |

Then denoting $r=\tau_{4}$ and $s=\tau_{2}$ we get $r^{4}=s^{2}=(s r)^{2}=1$ and

$$
\begin{aligned}
\overline{1} & =s^{2}=r^{4}, \\
\overline{\tau_{1}} & =r^{2}, \\
\overline{\tau_{2}} & =s, \\
\overline{\tau_{3}} & =s r, \\
\overline{\tau_{4}} & =r, \\
\overline{\tau_{1} \tau_{2}} & =r^{2} s, \\
\overline{\tau_{1} \tau_{3}} & =r^{2} s r=r s, \\
\overline{\tau_{1} \tau_{4}} & =r^{3} ;
\end{aligned}
$$

witnessing the desired isomorphism.
d-By the very definition

$$
\bigcap V=\left\{g=\left(t_{i}^{\alpha_{i}}\right)\left(t_{i j}^{\beta_{i j}}\right)\left(x_{i}^{\gamma_{i}}\right) \in \mathcal{W}(I): \beta_{i j}=\alpha_{i}=\gamma_{i}=\gamma_{j}=0 \text { for all } i, j \in I\right\}=\{1\} .
$$

Let $V$ as in Proposition 7.1.5 and let $P_{\text {fin }}(V)$ be the set of finite subsets of $V$ and for $A \in$ $P_{\text {fin }}(V)$, denote

$$
X_{A}=W_{I} / \bigcap A
$$

Note that $P_{\text {fin }}(V)$ is a directed poset with the partial ordering induced by inclusion. If $B \subseteq A$, denote by $\pi_{A B}: X_{A} \rightarrow X_{B}$ the canonical projection. Then $\left(X_{A}, \pi_{A B}, P_{f i n}(V)\right)$ is a projective system, in the sense that $\pi_{A A}=i d_{X_{A}}$ and, if $E \subseteq B \subseteq A$, then $\pi_{A E}=\pi_{B E} \circ \pi_{A B}$.
fixms3
Proposition 7.1.6. The canonical "diagonal" function

$$
\mathcal{W}(I) \rightarrow \lim _{A \in \stackrel{\leftrightarrow}{\text { fin }}(V)} X_{A} \text {, which is given by the rule } g \mapsto\left(g / X_{A}\right)_{A \in P_{\text {fin }}(V)}
$$

is an abstract group isomorphism, so, by transport, $\mathcal{W}(I)$ is a (topological) profinite 2-group with $P_{f i n}(V)$ as fundamental system of clopen neighborhoods of $\{1\}$.

Proof. Let $X=\prod_{A \in P_{\text {fin }}(V)} X_{A}$ and $\pi_{A}: X \rightarrow X_{A}$ be the canonical projection. Denote $\Delta: \mathcal{W}(I) \rightarrow$ $X$ the morphism given by the rule $\Delta(g):=\left(g / X_{A}\right)_{A \in P_{\text {fin }}(V)}$. This morphism $\Delta$ is injective since

$$
\operatorname{Ker}(\Delta)=\bigcap_{A \in P_{f i n}(V)}=\{1\}
$$

Now, let $\bar{g}=\left(g / X_{A}\right)_{A \in P_{f i n}(V)} \in \operatorname{Im}(\Delta)$. If $B \subseteq A$, we get

$$
\bar{g}_{B}=g / X_{B}=\left(g / X_{A}\right) / X_{B}=\left(\pi_{A B}(g)\right) / X_{B} .
$$

Moreover $\operatorname{Im}(\Delta) \subseteq \lim _{A \in P_{f i n}(V)} X_{A}$. To prove the surjectivity of $\Delta$, consider the morphism $\pi_{A}: \mathcal{W}(I) \rightarrow X_{A}$ given by the canonical projection. Then $\left(\mathcal{W}(I), \pi_{A}\right)$ is a compatible system of surjective morphisms where $\Delta$ is the exact morphism induced by $\left(\mathcal{W}(I), \pi_{A}\right)$. Then $\operatorname{Im}(\Delta)$ is a dense subset of ${\underset{\lim }{\rightleftarrows}}^{A \in P_{\text {fin }}(V)} X_{A}$ (for instance, see Lemma 1.1.7 of [56]) which is also closed. If $\varphi_{A}: \lim _{A \in P_{f i n}(V)} X_{A} \rightarrow X_{A}$ denote the projection, we have a new projective system $\left(\varphi_{A}\left(\operatorname{Im}(\Delta),\left.\pi_{A B}\right|_{\varphi_{A}(\operatorname{Im}(\Delta)}\right)\right.$. Then (using Corollary 1.1.8 of [56]) we get

$$
\operatorname{Im}(\Delta)=\lim _{A \in{\underset{P}{\text { fin }}}(V)} \varphi_{A}\left(\operatorname{Im}(\Delta)=\overline{\operatorname{Im}(\Delta)}=\varliminf_{A \in{\underset{P}{\text { fin }}}(V)}^{\lim _{A}} X_{A}\right.
$$

Moreover

$$
\left\{\operatorname{Ker}\left(\pi_{A}\right): \mathcal{W}(I) \rightarrow X_{A}\right\}_{A \in P_{f i n}(V)}=P_{f i n}(V)
$$

is a fundamental system of neighborhoods of $\{1\}$.
Lets invoke some terminology from the theory of profinite groups:
Definition 7.1.7. Let $G$ be a profinite group.
$i$ - We say that $X$ generates $G$ as a profinite group if $G=\overline{\langle G\rangle}$. In that case, we call $X a$ set of topological generators of $G$.
ii - We say that $X \subseteq G$ converges to 1 if every open subgroup $U$ of $G$ contains all but a finite number of the elements in $X$.
iii - Let $G$ be a profinite group. The Frattini subgroup of $G$, notation $\Phi(G)$, is the intersection of all its maximal open subgroups.

Fact 7.1.8. Let $G$ be a profinite group.
$i-G$ is compact, Hausdorff and totally disconnected (has a topological basis of clopens).
ii - $A$ subgroup $U$ is open if and only if is closed of finite index (Lemma 2.1.2 of [56]).
iii - A closed subgroup $H$ of a profinite group $G$ is the intersection of all open subgroups of $G$ containing $H$. If $H$ is normal, then $H$ is the intersection of all open normal subgroups of $G$ containing $H$ (Proposition 2.1.4 of [56]).
iv - A maximal closed subgroup is necessarily open.
$v-\Phi(G)$ is a characteristic subgroup of $G$ : for every automorphism $\psi: G \rightarrow G$ of $G$ we have $\psi[\Phi(G)]=\Phi(G)$.
vi- If $h: H \rightarrow G$ is a continuous homomorphism of pro-2-groups then $h[\Phi(H)] \subseteq \Phi[G]$.
rz2.8.7
Lemma 7.1.9 (Lemma 2.8.7 of [56]). Let $p$ be a prime number and let $G$ be a pro-p group.
a-Every maximal closed subgroup $M$ of $G$ has index $p$.
$b$ - The Frattini quotient $G / \Phi(G)$ is a p-elementary abelian profinite group and hence a vector space of the field $\mathbb{F}_{p}$ with $p$ elements.
$c-\Phi(G)=\overline{G^{p}[G, G]}$, where $G^{p}=\left\{x^{p}: x \in G\right\}$ and $[G, G]$ denotes the commutator subgroup of $G$.
lifting-le
Lemma 7.1.10. Let $\mathcal{G}_{i}, i=0,1$ be projectives profinite groups and $\mathcal{V}_{i} \subseteq \mathcal{G}_{i}$ be normal closed subgroups such that $\mathcal{V}_{i} \subseteq \Phi\left(\mathcal{G}_{i}\right)$. If $f: \mathcal{G}_{0} / \mathcal{V}_{0} \rightarrow \mathcal{G}_{1} / \mathcal{V}_{1}$ is an epimorphism (respectively an isomorphism) then there is some continuous homomorphism $\tilde{f}: \mathcal{G}_{0} \rightarrow \mathcal{G}_{1}$ such that $q_{1} \circ \tilde{f}=f \circ q_{0}$ where the $q_{i}$ are the projections on quotient; besides any such lifting $\tilde{f}$ is an epimorphism (resp. an isomorphism).


In [53] is established the main categorical property of $\mathcal{W}(I)$ :
Theorem 7.1.11 (Universal Property of $\mathcal{W}(I)$, Theorem 1.1 of [53]). The group $\mathcal{W}(I)$ is the $\mathcal{C}$ free group on I-generators. In other words, $I \subseteq \mathcal{W}(I)$ is a generator set converging to 1 and if $f: I \rightarrow G$ is any function to a $\mathcal{C}$-group $G$ such that $f[I] \subseteq G$ converges to 1 then there is an (unique) continuous homomorphism $\tilde{f}: \mathcal{W}(I) \rightarrow G$ such that $\left.\tilde{f}\right|_{I}=f$. Moreover, if $H$ is any $\mathcal{C}$-group then $H$ has a generator set converging to 1 of cardinality $|I|$ if and only if there is an epimorphism $\mathcal{W}(I) \rightarrow H$ with kernel $V$ contained in $\Phi(I)$, the Frattini subgroup of $\mathcal{W}(I)$.

Corollary 7.1.12. Let $X=\left\{x_{i}: i \in I\right\} \subseteq \mathcal{W}(I)$. Then $X$ is a set of generators of $\mathcal{W}(I)$ converging to 1 .
fixms5
Proposition 7.1.13. We have $\Phi(I)=\mathcal{W}_{I}^{2}$, and $\Phi(I)$ has $\left\{x_{i}, t_{i j}: i<j \in I\right\}$ as a minimal set of generators converging to 1. Moreover

$$
\begin{aligned}
\Phi(I) & =\mathcal{W}_{I}^{2}=\bigcap M(I) \\
& =\left\{g=\left(t_{i}^{\alpha_{i}}\right)\left(t_{i j}^{\beta_{i j}}\right)\left(x_{i}^{\gamma_{i}}\right) \in \mathcal{W}(I): \gamma_{i}=0 \text { for all } i \in I\right\} \\
& =\left\{g \in \mathcal{W}(I): \text { there exists } g_{1}, g_{2}, g_{3} \in \mathcal{W}(I) \text { such that } g=g_{1}^{2} g_{2}^{2} g_{3}^{2}\right\} .
\end{aligned}
$$

We have some kind of duality theorems for $\Phi(I)$ and $\mathcal{W}(I)$.
fixms6
Proposition 7.1.14. Consider the $\mathbb{Z}_{2}$-module of homogeneous quadratic polynomials in I variables $\left\{z_{i}\right\}_{i \in I}$

$$
P_{2}(I)=\left\{q \in \mathbb{Z}_{2}[I]: q=\sum_{i \in I} a_{i} z_{i}^{2}+\sum_{i<j \in I} b_{i j} z_{i} z_{j}\right\} \cong \bigoplus_{i \in I} \mathbb{Z}_{2} \oplus \bigoplus_{i<j \in I} \mathbb{Z}_{2} .
$$

Then we have a topological group isomorphism

$$
\Phi(I) \cong \operatorname{Hom}\left(P_{2}(I), \mathbb{Z}_{2}\right) \cong \prod_{i \in I} \mathbb{Z}_{2} \times \prod_{i<j \in I} \mathbb{Z}_{2},
$$

with the associated "perfect pairing"

$$
\langle-,-\rangle: \Phi(I) \times P_{2}(I) \rightarrow \mathbb{Z}_{2}
$$

Proof. First of all, note that $P_{2}(I)$ is generated (as $\mathbb{Z}_{2}$-vector space) by the set of monomials $B=\left\{z_{i}^{2}, z_{i} z_{j}\right\}_{i<j \in I}$. In fact, this is a $\mathbb{Z}_{2}$-basis of $P_{2}(I)$. Let $B^{*}$ be the dual basis

$$
B^{*}:=\left\{\varphi_{i j}\right\}_{i \leq j \in I},
$$

where $\varphi_{i j}: P_{2}(I) \rightarrow \mathbb{Z}_{2}$ is given by

$$
\varphi_{i j}\left(z_{k} z_{l}\right)=\left\{\begin{array}{l}
1 \text { if } i=k \text { and } j=l \\
0 \text { otherwise }
\end{array}\right.
$$

Now we define a function $\lambda:\left\{t_{i}, t_{i j}: i<j \in I\right\} \rightarrow B^{*}$ by the rules $t_{i} \mapsto \varphi_{i i}$ and $t_{i j} \mapsto \varphi_{i j}$. Since $\left\{t_{i}, t_{i j}: i<j \in I\right\}$ is a set of generators of $\Phi(I)$, this function $\lambda$ induces a continuous group homomorphism $\tilde{\lambda}: \Phi(I) \rightarrow \operatorname{Hom}\left(P_{2}(I), \mathbb{Z}_{2}\right)$ by the following: let $g=\left(t_{i}^{\alpha_{i}}\right)\left(t_{i j}^{\beta_{i j}}\right)\left(1_{i}\right) \in \Phi(I)$ (see Proposition 7.1.13. Define $\tilde{\lambda}(g): P_{2}(I) \rightarrow \mathbb{Z}_{2}$ for $q=\sum_{i \in I} a_{i} z_{i}^{2}+\sum_{i<j \in I} b_{i j} z_{i} z_{j} \in P_{2}(I)$ by

$$
\tilde{\lambda}(g)(q):=\sum_{i \in I} a_{i} \alpha_{i}+\sum_{i<j \in I} b_{i j} \beta_{i j} .
$$

We immediately have that $\tilde{\lambda}$ is a continuous injective group homomorphism. Since $\lambda$ is bijective and $B^{*}$ is a $\mathbb{Z}_{2}$-basis of $\operatorname{Hom}\left(P_{2}(I), \mathbb{Z}_{2}\right)$, we have that $\tilde{\lambda}$ is an isomorphism.

Proposition 7.1.15. Consider the $\mathbb{Z}_{2}$-module of homogeneous linear polynomials in I variables
$\left\{z_{i}\right\}_{i \in I}$

$$
P_{1}(I)=\left\{q \in \mathbb{Z}_{2}[I]: q=\sum_{i \in I} c_{i} z_{i}\right\} \cong \bigoplus_{i \in I} \mathbb{Z}_{2} .
$$

Then we have a topological group isomorphism

$$
\mathcal{W}(I) / \Phi(I) \cong \operatorname{Hom}\left(P_{1}(I), \mathbb{Z}_{2}\right) \cong \prod_{i \in I} \mathbb{Z}_{2},
$$

with the associated "perfect pairing"

$$
\langle-,-\rangle: \mathcal{W}(I) / \Phi(I) \times P_{1}(I) \rightarrow \mathbb{Z}_{2} .
$$

Proof. First of all, note that $P_{1}(I)$ is generated (as $\mathbb{Z}_{2}$-vector space) by the set of monomials $B=\left\{z_{i}\right\}_{i \in I}$. In fact, this is a $\mathbb{Z}_{2}$-basis of $P_{1}(I)$. Let $B^{*}$ be the dual basis

$$
B^{*}:=\left\{\varphi_{i}\right\}_{i \in I},
$$

where $\varphi_{i}: P_{1}(I) \rightarrow \mathbb{Z}_{2}$ is given by

$$
\varphi_{i}\left(z_{j}\right)=\left\{\begin{array}{l}
1 \text { if } i=j \\
0 \text { otherwise }
\end{array}\right.
$$

Now we define $\theta: \mathcal{W}(I) \rightarrow \operatorname{Hom}\left(P_{1}(I), \mathbb{Z}_{2}\right)$ by the following: for $g=\left(t_{i}^{\alpha_{i}}\right)\left(t_{i j}^{\beta_{i j}}\right)\left(x_{i}^{\gamma_{i}}\right) \in \mathcal{W}(I)$, let $\theta(g): P_{1}(I) \rightarrow \mathbb{Z}_{2}$ be the morphism defined by the rule

$$
\text { For } q=\sum_{i \in I} c_{i} z_{i} \in P_{1}(I), \theta(g)(q)=\sum_{i \in I} c_{i} \gamma_{i} \text {. }
$$

Then $\theta$ is a surjective morphism (because $B^{*} \subseteq \operatorname{Im}(\theta)$ ) and by Proposition 7.1.13 $\operatorname{Ker}(\theta)=\Phi(I)$. Hence

$$
\mathcal{W}(I) / \Phi(I)=\mathcal{W}(I) / \operatorname{Ker}(\theta) \cong \operatorname{Hom}\left(P_{1}(I), \mathbb{Z}_{2}\right) .
$$

Now, is time to return to our first goal: present the description of $G_{F}:=\operatorname{Gal}_{F}\left(F^{(3)}\right)$ by $G_{F} \cong \mathcal{W}(I) / \mathcal{V}(I)$.

Proposition 7.1.16 (2.1 in 53]). $G_{F} \cong \overline{G_{F}^{q}}$.
Then $G_{F}$ is a $C$-group on $B$, where $B=\left\{a_{i}: i \in I\right\}$ is an well-ordered basis of $\dot{F} / \dot{F}^{2}$, so, using Theorem 7.1.11, there exists an epimorphism $\pi_{B}: \mathcal{W}(B) \rightarrow G_{F}$. Then we simply take $\mathcal{V}(B):=\operatorname{Ker}\left(\pi_{B}\right)$.

Moreover, J. Minac and M. Spira (again, in [53]) gave a nice explicit description to $\mathcal{V}(B)$. Let Quat $(F)$ be the subgroup of $\operatorname{Br}(F)$, the Brauer group of $F$, generated by the quaternion algebras of $F$. By Merkurjev's Theorem (59]), we have

$$
\operatorname{Quat}(F) \cong \operatorname{Br}_{2}(F) \cong k_{2}(F),
$$

where $B r_{2}(F)$ is the subgroup of $\operatorname{Br}(F)$ generated by elements of order 2 and $k_{*}(F)$ is the graded ring of Milnor's mod 2 reduced K-theory.

Consider $\varphi_{B}: P_{2}(B) \rightarrow k_{2}(F)$ as the epimorphism defined by the rule

$$
\left(\sum_{i \in I} \alpha_{i} z_{i}^{2}+\sum_{i<j \in I} \beta_{i j} z_{i} z_{j}\right) \mapsto\left(\sum_{i \in I} \alpha_{i} l\left(a_{i}\right) l\left(a_{i}\right)+\sum_{i<j \in I} \beta_{i j} l\left(a_{i}\right) l\left(a_{j}\right)\right) .
$$

Let $Q_{B}:=\operatorname{Ker}\left(\varphi_{B}\right)$.
Fact 7.1.17 (Essentially 2.20 in [53]). $\mathcal{V}(B)=Q_{B}^{\perp}$, where $Q_{B}^{\perp}=\left\{v \in \Phi(B):\left\langle v, Q_{B}\right\rangle=0\right\}$ and $\langle$,$\rangle is the perfect pairing described in Proposition 7.1.14.$

Let a $F$ be a field with $\operatorname{char}(F) \neq 2$. By "Pontryaguin duality", let $M_{F}$ denotes the unique maximal clopen subgroup of $F_{F}=\operatorname{Gal}_{F}\left(F^{(3)}\right)$ corresponding to $-1 \in S G(F)=\dot{F} / \dot{F}^{2}$ (Proposition 7.1.15).

Fact 7.1.18 (Essentially 3.3, 3.5, 3.6, 3.7 in [53]). Let $F, L \in$ Field $_{2}$. Then are equivalent:
$a-\left(W(F), \dot{F} / \dot{F}^{2}\right) \cong\left(W(L), \dot{L} / \dot{L}^{2}\right)$ as abstract Witt rings.
$b-S G(F) \cong S G(L)$ as special groups.
$c-\left(G_{F}, M_{F}\right) \cong\left(G_{L}, M_{L}\right)$ as pointed profinite- $\mathcal{C}$-groups.
Our next step, is use all these facts to obtain a group associated to a (pre)-special group.

### 7.2 The Galois Group of a Pre Special Group

Lets deal first, with a special group $G$. Let $B=\left\{a_{i}: i \in I\right\}$ be a well ordered $\mathbb{Z}_{2}$-basis of $G$ and consider the $\mathcal{C}$-free group in $B$-generators $\mathcal{W}(B)$. Define an epimorphism $\pi_{B}: P_{2}(B) \rightarrow k_{2}(G)$ by the rule

$$
\left(\sum_{i \in I} \alpha_{i} z_{i}^{2}+\sum_{i<j \in I} \beta_{i j} z_{i} z_{j}\right) \mapsto\left(\sum_{i \in I} \alpha_{i} l\left(a_{i}\right) l\left(a_{i}\right)+\sum_{i<j \in I} \beta_{i j} l\left(a_{i}\right) l\left(a_{j}\right)\right)
$$

with kernel $Q(B)$. Take $\mathcal{V}(B):=Q(B)^{\perp} \subseteq \Phi(B) \subseteq \mathcal{W}(B)$. Since $\Phi(B)$ is the center of $\mathcal{W}(B)$ then $\mathcal{V}(B) \subseteq \mathcal{W}(B)$ is a (closed) normal subgroup of $\mathcal{W}(B)$ and we can consider the $\mathcal{C}$-group $\mathcal{W}(B) / \mathcal{V}(B)$.
defn:gal1
Definition 7.2.1 (Galois Group - base dependent version). Let $G$ be a special group and $B, \mathcal{W}(B)$ and $\mathcal{V}(B)$ as above. We define the Galois group of $G$ with respect to $B$ by

$$
\operatorname{Gal}(G, B):=\mathcal{W}(B) / \mathcal{V}(B)
$$

The most essential information of our Galois group is encoded by $Q(B)=\operatorname{Ker}\left(\pi_{B}\right)$. We have a useful description by generators that generalizes the one described by J. Minac and M. Spira:

Proposition 7.2.2. Let $G$ be a special group and $B, \mathcal{W}(B)$ and $\mathcal{V}(B)$ as above. Consider a finite subset $B^{\prime} \subseteq B, B^{\prime}=\left\{a_{i_{0}}, \ldots, a_{i_{n-1}}\right\}\left(i_{0}<\ldots<i_{n-1}\right)$, and $a, b$ in the linear span of $B^{\prime}$, say

$$
a=\prod_{k<n} a_{i_{k}}^{\alpha_{i k}} \text { and } b=\prod_{k<n} a_{i_{k}}^{\beta_{i_{k}}}, \alpha_{i_{k}}, \beta_{i_{k}} \in\{0,1\} .
$$

Consider the polynomial $q_{a, b}^{B} \in P_{2}(B)$ given by

$$
q_{a, b}^{B}=\sum_{k<n} \alpha_{i_{k}} \beta_{i_{k}} z_{i_{k}}^{2}+\sum_{k<l<n}\left(\alpha_{i_{k}} \beta_{i_{l}}+\alpha_{i_{l}} \beta_{i_{k}}\right) z_{i_{k}} z_{i_{l}} .
$$

Note that $q_{a,\left(b_{0} \ldots b_{n-1}\right)}^{B}=\sum_{k<n} q_{a, b_{k}}^{B}$. Moreover we have the following properties.
$i-\pi_{B}\left(q_{a, b}^{B}\right)=l(a) l(b) \in k_{2}(G)$.
ii - $q_{a, b}^{B}$ does not depend on the particular choice of the finite subset $B^{\prime} \subseteq B$.
iii $-Q(B)$ is generated by $\left\{q_{a, b}^{B}: l(a) l(b)=0\right\}$.
Proof.
i - Note that

$$
l(a)=\sum_{k<n} \alpha_{i_{k}} l\left(a_{i_{k}}\right) \text { and } l(b)=\sum_{k<n} \beta_{i_{k}} l\left(a_{i_{k}}\right) .
$$

Then

$$
\begin{aligned}
l(a) l(b) & =\left(\sum_{k<n} \alpha_{i_{k}} l\left(a_{i_{k}}\right)\right)\left(\sum_{k<n} \beta_{i_{k}} l\left(a_{i_{k}}\right)\right)=\sum_{k<n} \sum_{p<n} \alpha_{i_{k}} \beta_{i_{p}} l\left(a_{i_{k}}\right) l\left(a_{i_{p}}\right) \\
& =\sum_{k<n} \alpha_{i_{k}} \beta_{i_{k}} l\left(a_{i_{k}}\right) l\left(a_{i_{k}}\right)+\sum_{k<n} \sum_{k<p<n} \alpha_{i_{k}} \beta_{i_{p}} l\left(a_{i_{k}}\right) l\left(a_{i_{p}}\right)+\sum_{k<n} \sum_{p<k<n} \alpha_{i_{k}} \beta_{i_{p}} l\left(a_{i_{k}}\right) l\left(a_{i_{p}}\right) \\
& =\sum_{k<n} \alpha_{i_{k}} \beta_{i_{k}} l\left(a_{i_{k}}\right) l\left(a_{i_{k}}\right)+\sum_{k<n} \sum_{k<p<n}\left(\alpha_{i_{k}} \beta_{i_{p}}+\alpha_{i_{p}} \beta_{i_{k}}\right) l\left(a_{i_{k}}\right) l\left(a_{i_{p}}\right)
\end{aligned}
$$

On the other hand, by definition of $\pi_{B}$ we get

$$
\pi_{B}\left(q_{a, b}^{B}\right)=\sum_{k<n} \alpha_{i_{k}} \beta_{i_{k}} l\left(a_{i_{k}}\right) l\left(a_{i_{k}}\right)+\sum_{k<p<n}\left(\alpha_{i_{k}} \beta_{i_{p}}+\alpha_{i_{p}} \beta_{i_{k}}\right) l\left(a_{i_{k}}\right) l\left(a_{i_{p}}\right),
$$

completing the proof.
ii - It is an immediate consequence of previous item: if $B_{1}, B_{2}$ are finite subsets of $B$ and $a, b$ are elements in the linear span of $B_{1}$ and in the linear span of $B_{2}$, then

$$
q_{a, b}^{B_{1}}=q_{a, b}^{B_{2}} .
$$

iii - Of course, $q_{a, b}^{B} \in Q(B)$ if and only if $l(a) l(b)=0$ in $k_{2}(G)$ and hence

$$
\left[\left\{q_{a, b}^{B}: l(a) l(b)=0\right\}\right] \subseteq Q(B)
$$

To get the reverse inclusion, let $q=\sum_{k<n} \alpha_{i_{k}} z_{i_{k}}^{2}+\sum_{k, p<n} \beta_{i_{k} i_{p}} z_{i_{k}} z_{i_{p}} \in Q$. Then

$$
\sum_{k<n} \alpha_{i_{k}} l\left(a_{i_{k}}\right) l\left(a_{i_{k}}\right)+\sum_{k, p<n} \beta_{i_{k} i_{p}} l\left(a_{i_{k}}\right) l\left(a_{i_{p}}\right)=0 \text { in } k_{2}(G) .
$$

Now, for each $k<n$ let

$$
b_{k}:=a_{i_{k}}^{\alpha_{i_{k}}} \prod_{k<p} a_{i_{p}}^{\beta_{i_{k} i_{p}}} .
$$

Then $q=\sum_{k<n} q_{a_{i_{k}}, b_{i_{k}}}^{B}$ and

$$
\sum_{k<n} l\left(a_{i_{k}}\right) l\left(b_{i_{k}}\right)=0 \text { in } k_{2}(G) .
$$

We are under the hypothesis of Lemma 7.2.3. Thus, according Theorem 4.3.8, there exists subsets $\left\{c_{0}, \ldots, c_{m-1}\right\},\left\{d_{0}, \ldots, d_{n-1}\right\}$ of $G$ with $m \geq n$ such that
(a) $\left\{c_{0}, \ldots, c_{m-1}\right\}$ is linearly independent and $c_{k}=a_{i_{k}}$ for all $k<n$;
(b) $d_{k}=b_{i_{k}}$ for all $k<n$ and $d_{k}=1$ for $k=n, \ldots, m-1$.
(c) For all $x \in\left[c_{0}, \ldots, c_{m-1}\right]$, there is some $r_{x} \in D_{G}(1,-x)$ such that for each $k<m$

$$
d_{k}=\prod_{x \in C_{k}} r_{x}
$$

where

$$
C_{k}=\left\{\prod_{p<m} c_{p}^{\varepsilon_{p}}: \varepsilon_{p} \in\{0,1\} \text { and } \varepsilon_{k}=1\right\} .
$$

It follows that

$$
\begin{aligned}
q & =\sum_{k<n} q_{a_{i_{k}}, b_{i_{k}}}^{B}=\sum_{k<m} q_{c_{k}, d_{k}}^{B}=\sum_{k<m} q_{c_{k}, \Pi_{x \in C_{k}} r_{x}}^{B} \\
& =\sum_{k<m} \sum_{x \in C_{k}} q_{c_{k}, r_{x}}^{B} .
\end{aligned}
$$

Denoting $C:=\left[c_{0}, \ldots, c_{m-1}\right]$, we have $C=C_{0} \cup \ldots \cup C_{m-1}$. Then

$$
q=\sum_{k<m} \sum_{x \in C_{k}} q_{c_{k}, r_{x}}^{B}=\sum_{x \in C} q_{x, r_{x}}^{B} .
$$

Since $r_{x} \in D_{G}(1,-x)$, we have $l(x) l\left(r_{x}\right)=0$ in $k_{2}(G)$. Then

$$
q=\sum_{x \in C} q_{x, r_{x}}^{B} \in\left[\left\{q_{a, b}^{B}: l(a) l(b)=0\right\}\right] .
$$

Now, we will generalize the Galois group for pre-special groups. The K-theory developed by M. Dickmann and F. Miraglia in [30] is available for pre-special groups. Then we can take the same $B, \mathcal{W}(B)$ and $\mathcal{V}(B)$ we are considering until now.

Let $G$ be a special group and $B=\left\{v_{i}\right\}_{i \in I}, C=\left\{w_{i}\right\}_{i \in I}, D=\left\{z_{i}\right\}_{i \in I}$ be ordered $\mathbb{Z}_{2}$-basis of $G$. Then, for all $i \in I$ we have an expression

$$
w_{i}=\prod_{k \in I} v_{k}^{m_{i k}}, m_{i k} \in\{0,1\} \text { for all } i, k \in I
$$

base-change-prod
where the above product has finite support (i.e, $\left|\left\{i, k \in I: m_{i k} \neq 0\right\}\right|<\infty$ ). In other words, for all $i \in I$, there exist unique sequence in $I i_{0}<i_{1}<\cdots<i_{n}$ such that

$$
w_{i}=v_{i_{0}} \cdot v_{i_{1}} \cdot \ldots \cdot v_{i_{n(i)}} .
$$

By abuse of notation, let $C=\left\{x_{i}: i \in I\right\} \subseteq \mathcal{W}(C)$. We define a function $\mu_{B C}: C \rightarrow \mathcal{W}(B)$ by the rule

$$
x_{i} \mapsto x_{i_{0}} \cdot x_{i_{1}} \cdot \ldots \cdot x_{i_{n}} \text { if } w_{i}=v_{i_{0}} \cdot v_{i_{1}} \cdot \ldots \cdot v_{i_{n}} .
$$

This function is well-defined because both $B$ and $C$ are basis, so the expression 7.2 is unique.
fixsg2
Lemma 7.2.3. Let $G$ be a pre-special group and $B, C, \mu_{B C}$ as above. There is a unique continuous homomorphism $\mu_{B C}: \mathcal{W}(C) \rightarrow \mathcal{W}(B)$ that extends $\mu_{B C}$. Also $\mu_{B C}[\Phi(C)] \subseteq \Phi(B)$.

Proof. By abuse of notation, let $B=\left\{x_{i}: i \in I\right\} \subseteq \mathcal{W}(B)$ and $C=\left\{x_{i}: i \in I\right\} \subseteq \mathcal{W}(C)$. We have $B$ and $C$ as a set of generators converging to 1 in $\mathcal{W}(B)$ and $\mathcal{W}(C)$ respectively.

Let $X=\mu_{B C}[C] \subseteq \mathcal{W}(B)$. Since $\langle X\rangle=\langle B\rangle$ and $\mathcal{W}(B)=\overline{\langle B\rangle}$, we have that $X$ is a set of generators of $\mathcal{W}(B)$.

Let $U \subseteq \mathcal{W}(B)$ be an open subgroup. Since $B$ is a set of generators converging to 1 , there is a finite subset $Y=\left\{x_{i_{1}}, \ldots, x_{j_{m}}\right\} \subseteq B$ with $U \cap Y=\emptyset$. Since the set of $\mathbb{F}_{2}$-linear combinations of a finite set is finite, there is a finite quantity of elements in $\mu_{B C}[C]$ not belonging to $U$.

The existence and continuity of $\mu_{B C}$ is an immediate consequence of the Universal Property of $\mathcal{W}(I)$ (Theorem 7.1.11). An explicit formula for $\mu_{B C}$ is given by the following rule: for $g=$ $\left(t_{i}^{\alpha_{i}}\right)\left(t_{i j}^{\beta_{i j}}\right)\left(x_{i}^{\gamma_{i}}\right)$ we set

$$
\mu_{B C}(g):=\left(t_{i}^{\left(\alpha_{i} \sum_{k \in I} m_{i k}\right)}\right)\left(t_{i j}^{\left(\beta_{i j} \sum_{r, s \in I} m_{i r} m_{j s}\right)}\right)\left(x_{i}^{\left(\gamma_{i} \sum_{k \in I} m_{i k}\right)}\right) .
$$

Since $\Phi(\mathcal{W}(C))=\mathcal{W}(C)^{2}$ and $\Phi(\mathcal{W}(B))=\mathcal{W}(B)^{2}$, we get $\mu_{B C}[\Phi(C)] \subseteq \Phi(B)$.

Remark 7.2.4. A direct calculation show for all $a, b \in G$ that

$$
m_{B, B^{\prime}}^{2}\left(q_{a b}^{B}\right)=q_{a b}^{B^{\prime}} .
$$

Denote $\mu_{B, B^{\prime}}^{(2)}: \Phi\left(B^{\prime}\right) \rightarrow \Phi(B)$ the restriction of $\mu_{B, B^{\prime}}$ to the Frattini's subgroups and $\mu_{B, B^{\prime}}^{(1)}$ : $\mathcal{W}\left(B^{\prime}\right) / \Phi\left(B^{\prime}\right) \rightarrow \mathcal{W}(B) / \Phi(B)$ the quotient of $\mu_{B, B^{\prime}}$. Then, from the isomorphism in Proposition 7.1.14 $\Phi(B) \cong \operatorname{Hom}\left(P_{2}(B), \mathbb{Z}_{2}\right)$, we have:

$$
\Phi\left(B^{\prime}\right) \times P_{2}(B) \rightarrow \mathbb{Z}_{2}:<\mu_{B, B^{\prime}}^{1}-,->_{B}=<-, m_{B, B^{\prime}}^{1}->_{B^{\prime}}
$$

so, for all $\sigma^{\prime} \in \Phi\left(B^{\prime}\right)$

$$
<\mu_{B, B^{\prime}}^{2}\left(\sigma^{\prime}\right), q_{a b}^{B}>_{B}=<\sigma^{\prime}, m_{B, B^{\prime}}^{2}\left(q_{a b}^{B}\right)>_{B^{\prime}}=<\sigma^{\prime}, q_{a b}^{B^{\prime}}>_{B^{\prime}}
$$

and then $\mu_{B, B^{\prime}}^{2}\left[\left(q_{a, b}^{B^{\prime}}\right)^{\perp}\right]=\left(q_{a b}^{B}\right)^{\perp}$

Lemma 7.2.5. The morphism $\mu_{B C}$ is an isomorphism, $\mu_{B B}=i d_{\mathcal{W}(B)}, \mu_{B C}^{-1}=\mu_{C B}$ and

$$
\mu_{B D}=\mu_{B C} \circ \mu_{C D} .
$$

Proof. The fact that $\mu_{B B}=i d_{\mathcal{W}(B)}$ and $\mu_{B C}^{-1}=\mu_{C B}$ is an immediate consequence of Lemma 7.2.3. For the other part, let $B=\left\{v_{i}\right\}_{i \in I}, C=\left\{w_{i}\right\}_{i \in I}$ and $D=\left\{z_{i}\right\}_{i \in I}$ be $\mathbb{F}_{2}$-basis of $G$. Then for all
$i \in I$,

$$
w_{i}=\prod_{k \in I} v_{k}^{m_{i k}}, z_{i}=\prod_{k \in I} w_{k}^{n_{i k}}, z_{i}=\prod_{k \in I} v_{k}^{p_{i k}},
$$

such that all these products has finite support. Then

$$
z_{i}=\prod_{k \in I} w_{k}^{n_{i k}}=\prod_{k \in I}\left(\prod_{r \in I} v_{r}^{m_{k r}}\right)^{n_{i k}}=\prod_{k \in I} \prod_{r \in I} v_{r}^{n_{i k} m_{k r}}=\prod_{r \in I} \prod_{k \in I} v_{r}^{n_{i k} m_{k r}}=\prod_{r \in I} v_{r}^{p_{i r}} .
$$

Moreover

$$
\sum_{k \in I} \sum_{r \in I} n_{i k} m_{k r}=\sum_{r \in I} \sum_{k \in I} n_{i k} m_{k r}=\sum_{r \in I} p_{i r} .
$$

Then for all $g=\left(t_{i}^{\alpha_{i}}\right)\left(t_{i j}^{\beta_{i j}}\right)\left(x_{i}^{\gamma_{i}}\right) \in \mathcal{W}(D)$,

$$
\left.\begin{array}{l}
\mu_{B C} \circ \mu_{C D}(g)=\mu_{B C}\left(\mu_{C D}(g)\right)= \\
\mu_{B C}\left(\left(t_{i}^{\alpha_{i}\left(\sum_{k \in I} n_{i k}\right)}\right)\left(t_{i j}^{\left(\beta_{i j} \sum_{r, s \in I} n_{i r} n_{j s}\right)}\right)\left(x_{i}^{\left(\gamma_{i} \sum_{k \in I} n_{i k}\right)}\right)\right)= \\
\left(t_{i}^{\alpha_{i}\left(\sum_{k \in I} n_{i k}\right)\left(\sum_{r \in I} m_{k r}\right)}\right)\left(t_{i j}^{\beta_{i j}} \sum_{r, s \in I} n_{i r} n_{j s}\right)\left(\sum_{e, f \in I} m_{r e} m_{s f}\right)
\end{array}\right)\left(x_{i}^{\gamma_{i}\left(\sum_{k \in I} n_{i k}\right)\left(\sum_{r \in I} m_{k r}\right)}\right)=\begin{aligned}
& \left(t_{i}^{\alpha_{i}\left(\sum_{r \in I} \sum_{k \in I} n_{i k} m_{k r}\right)}\right)\left(t_{i j}^{\beta_{i j}\left(\sum_{r, s \in I} \sum_{e, f \in I}\left(n_{i r} m_{r e}\right)\left(n_{j s} m_{s f}\right)\right)}\right)\left(x_{i}^{\gamma_{i}\left(\sum_{r \in I} \sum_{k \in I} n_{i k} m_{k r}\right)}\right)= \\
& \left(t_{i}^{\alpha_{i}\left(\sum_{r \in I} p_{i r}\right)}\right)\left(t_{i j}^{\beta_{i j}\left(\sum_{e, f \in I} p_{i e} p_{j f}\right)}\right)\left(x_{i}^{\gamma_{i}\left(\sum_{r \in I} p_{i r}\right)}\right)=\mu_{B D}(g) .
\end{aligned}
$$

Then we get $\mu_{B D}=\mu_{B C} \circ \mu_{C D}$.

Now, consider the Equation 7.1. This expression induce isomorphisms $m_{B C}^{1}: P_{1}(B) \rightarrow P_{1}(C)$ and $m_{B C}^{2}: P_{2}(B) \rightarrow P_{2}(C)$ given by the rules

$$
\begin{aligned}
m_{B C}^{1}\left(\sum_{i \in I} c_{i} z_{i}\right) & :=\sum_{i \in I} c_{i}\left(\sum_{k \in I} m_{i k}\right) z_{i} \\
m_{B C}^{2}\left(\sum_{i \in I} a_{i} z_{i}^{2}+\sum_{i<j \in I} b_{i j} z_{i} z_{j}\right) & :=\sum_{i \in I} a_{i}\left(\sum_{k \in I} m_{i k}\right) z_{i}^{2}+\sum_{i<j \in I} b_{i j}\left[\left(\sum_{r \in I} m_{i r}\right)\left(\sum_{s \in I} m_{j s}\right)\right] z_{i} z_{j} .
\end{aligned}
$$

where all these sums has finite support.

Lemma 7.2.6. We have

$$
m_{B D}^{1}=m_{B C}^{1} \circ m_{C D}^{1} \text { and } m_{B D}^{2}=m_{B C}^{2} \circ m_{C D}^{2}
$$

Proof. Lets recover the calculations in the proof of Lemma 7.2.5. let $B=\left\{v_{i}\right\}_{i \in I}, C=\left\{w_{i}\right\}_{i \in I}$ and $D=\left\{z_{i}\right\}_{i \in I}$ be $\mathbb{F}_{2}$-basis of $G$. Then for all $i \in I$,

$$
w_{i}=\prod_{k \in I} v_{k}^{m_{i k}}, z_{i}=\prod_{k \in I} w_{k}^{n_{i k}}, z_{i}=\prod_{k \in I} v_{k}^{p_{i k}},
$$

such that all these products has finite support. Then

$$
z_{i}=\prod_{k \in I} w_{k}^{n_{i k}}=\prod_{k \in I}\left(\prod_{r \in I} v_{r}^{m_{k r}}\right)^{n_{i k}}=\prod_{k \in I} \prod_{r \in I} v_{r}^{n_{i k} m_{k r}}=\prod_{r \in I} \prod_{k \in I} v_{r}^{n_{i k} m_{k r}}=\prod_{r \in I} v_{r}^{p_{i r}}
$$

Moreover

$$
\sum_{k \in I} \sum_{r \in I} n_{i k} m_{k r}=\sum_{r \in I} \sum_{k \in I} n_{i k} m_{k r}=\sum_{r \in I} p_{i r} .
$$

Then

$$
\begin{aligned}
m_{B C}^{1}\left(m_{C D}^{1}\left(\sum_{k \in I} c_{k} z_{i}\right)\right) & =m_{B C}^{1}\left(\sum_{i \in I} c_{i}\left(\sum_{k \in I} n_{i k}\right) z_{i}\right)=\sum_{i \in I} c_{i}\left(\sum_{k \in I} n_{i k}\right)\left(\sum_{r \in I} m_{k r}\right) z_{i} \\
& =\left(\sum_{r \in I} p_{i r}\right) z_{i}=m_{B D}^{1}\left(\sum_{k \in I} c_{k} z_{i}\right)
\end{aligned}
$$

and hence $m_{B D}^{1}=m_{B C}^{1} \circ m_{C D}^{1}$. In the same reasoning,

$$
\begin{aligned}
& m_{B C}^{2}\left(m_{C D}^{2}\left(\sum_{i \in I} a_{i} z_{i}^{2}+\sum_{i<j \in I} b_{i j} z_{i} z_{j}\right)\right)= \\
& m_{B C}^{2}\left(\sum_{i \in I} a_{i}\left(\sum_{k \in I} n_{i k}\right) z_{i}^{2}+\sum_{i<j \in I} b_{i j}\left[\left(\sum_{r \in I} n_{i r}\right)\left(\sum_{s \in I} n_{j s}\right)\right] z_{i} z_{j}\right)= \\
& \sum_{i \in I} a_{i}\left(m_{B C}^{2}\left(\left(\sum_{k \in I} n_{i k}\right) z_{i}^{2}\right)\right)+m_{B C}^{2}\left(\sum_{i<j \in I} b_{i j}\left[\left(\sum_{r \in I} n_{i r}\right)\left(\sum_{s \in I} n_{j s}\right)\right] z_{i} z_{j}\right)= \\
& \sum_{i \in I} a_{i}\left(\sum_{k \in I} n_{i k}\left(\sum_{r \in I} m_{k r}\right) z_{i}^{2}\right)+\sum_{i<j \in I} b_{i j}\left[\left(\sum_{r \in I} n_{i r}\right)\left(\sum_{s \in I} n_{j s}\right)\left(\sum_{e \in I} m_{r e}\right)\left(\sum_{f \in I} m_{s f}\right)\right] z_{i} z_{j}= \\
& \sum_{i \in I} a_{i}\left(\left(\sum_{r \in I} \sum_{k \in I} n_{i k} m_{k r}\right) z_{i}^{2}\right)+\sum_{i<j \in I} b_{i j}\left[\left(\sum_{e \in I} \sum_{r \in I} n_{i r} m_{r e}\right)\left(\sum_{f \in I} \sum_{s \in I} n_{j s} m_{s f}\right)\right] z_{i} z_{j}= \\
& \sum_{i \in I} a_{i}\left(\sum_{r \in I} p_{i r}\right) z_{i}^{2}+\sum_{i<j \in I} b_{i j}\left[\left(\sum_{e \in I} p_{i e}\right)\left(\sum_{f \in I} p_{j f}\right)\right] z_{i} z_{j}= \\
& m_{B D}^{2}\left(\sum_{i \in I} a_{i} z_{i}^{2}+\sum_{i<j \in I} b_{i j} z_{i} z_{j}\right)
\end{aligned}
$$

provide that $m_{B D}^{2}=m_{B C}^{2} \circ m_{C D}^{2}$.

Note that $m_{B C}^{1}, m_{B C}^{2}$ induces respectively the isomorphisms

$$
\begin{aligned}
& \mathfrak{m}_{B C}^{1}: \operatorname{Hom}\left(P_{1}(C), \mathbb{Z}_{2}\right) \rightarrow \operatorname{Hom}\left(P_{1}(B), \mathbb{Z}_{2}\right) \\
& \mathfrak{m}_{B C}^{2}: \operatorname{Hom}\left(P_{2}(C), \mathbb{Z}_{2}\right) \rightarrow \operatorname{Hom}\left(P_{2}(B), \mathbb{Z}_{2}\right) .
\end{aligned}
$$

given by the respective rules: if $f: P_{1}(C) \rightarrow \mathbb{Z}_{2}$ and $q=\sum_{i \in I} c_{i} z_{i} \in P_{1}(B)$ then

$$
\mathfrak{m}_{B C}^{1}(f)(q):=f\left(m_{B C}^{1}(q)\right)=f\left(\sum_{i \in I} c_{i}\left(\sum_{k \in I} m_{i k}\right) z_{i}\right) .
$$

In the same reasoning, if $f: P_{2}(C) \rightarrow \mathbb{Z}_{2}$ and $q=\sum_{i \in I} a_{i} z_{i}^{2}+\sum_{i<j \in I} b_{i j} z_{i} z_{j} \in P_{2}(B)$ then

$$
\mathfrak{m}_{B C}^{2}(f)(q):=f\left(m_{B C}^{2}(q)\right)=f\left(\sum_{i \in I} a_{i}\left(\sum_{k \in I} m_{i k}\right) z_{i}^{2}+\sum_{i<j \in I} b_{i j}\left[\left(\sum_{r \in I} m_{i r}\right)\left(\sum_{s \in I} m_{j s}\right)\right] z_{i} z_{j}\right)
$$

Now denote

$$
\mu_{B C}^{1}: \mathcal{W}(C) / \Phi(C) \rightarrow \mathcal{W}(B) / \Phi(B)
$$

the quotient of $\mu_{B C}$ and

$$
\mu_{B C}^{2}: \Phi(C) \rightarrow \Phi(B)
$$

the restriction of $\mu_{B C}$ to the Frattini's subgroups. Also consider the isomorphisms

$$
\begin{aligned}
& \tilde{\theta}: \mathcal{W}(I) / \Phi(I) \cong \\
& \tilde{\lambda}: \Phi(I) \xlongequal{\cong} \operatorname{Hom}\left(P_{1}(I), \mathbb{Z}_{2}\right) \\
& \operatorname{Hom}\left(P_{2}(I), \mathbb{Z}_{2}\right)
\end{aligned}
$$

of Propositions 7.1.14, and 7.1.15.

Lemma 7.2.7. Denote $\pi_{B}: \mathcal{W}(B) \rightarrow \mathcal{W}(B) / \Phi(B)$ the canonical projection, with the same for $\pi_{C}$. Then we have a commutative diagram

which induces a commutative diagram


Proof. Let $g=\left(t_{i}^{\alpha_{i}}\right)\left(t_{i j}^{\beta_{i j}}\right)\left(x_{i}^{\gamma_{i}}\right) \in \mathcal{W}(C)$ and $q=\sum_{i \in I} c_{i} z_{i} \in P_{1}(B)$. Then

$$
\begin{aligned}
\left(\mathfrak{m}_{B C}^{1} \circ \theta\right)(g)(q) & =\mathfrak{m}_{B C}^{1}(\theta(g)(q))=\theta(g)\left(m_{B C}^{1}(q)\right) \\
& =\theta(g)\left(\sum_{i \in I} c_{i}\left(\sum_{k \in I} m_{i k}\right) z_{i}\right)=\sum_{i \in I} c_{i}\left(\sum_{k \in I} m_{i k}\right) \gamma_{i} \\
& =\sum_{i \in I} \sum_{k \in I} c_{i} m_{i k} \gamma_{i} .
\end{aligned}
$$

On the other hand,

$$
\begin{aligned}
\left(\theta \circ \mu_{B C}\right)(g)(q) & =\theta\left(\mu_{B C}(g)\right)(q) \\
& =\theta\left(\left(t_{i}^{\alpha_{i}\left(\sum_{k \in I} m_{i k}\right)}\right)\left(t_{i j}^{\beta_{i j}\left(\sum_{r, s \in I} m_{i r} m_{j s}\right)}\right)\left(x_{i}^{\gamma_{i}\left(\sum_{k \in I} m_{i k}\right)}\right)\right)\left(\sum_{i \in I} c_{i} z_{i}\right) \\
& =\sum_{i \in I} c_{i}\left(\gamma_{i}\left(\sum_{k \in I} m_{i k}\right)\right)=\sum_{i \in I} \sum_{k \in I} c_{i} m_{i k} \gamma_{i}=\left(\mathfrak{m}_{B C}^{1} \circ \theta\right)(g)(q) .
\end{aligned}
$$

Then $\mathfrak{m}_{B C}^{1} \circ \theta=\theta \circ \mu_{B C}$. Since $\theta_{B}=\tilde{\theta}_{B} \circ \pi_{B}$ and $\theta_{C}=\tilde{\theta}_{C} \circ \pi_{C}$ we have the desired commutative diagram.

Lemma 7.2.8. We have a commutative diagram


Proof. Let $g=\left(t_{i}^{\alpha_{i}}\right)\left(t_{i j}^{\beta_{i j}}\right)\left(1_{i}\right) \in \Phi(C)$ and $\sum_{i \in I} a_{i} z_{i}^{2}+\sum_{i<j \in I} b_{i j} z_{i} z_{j} \in P_{2}(B)$. We have

$$
\begin{aligned}
& \left(\mathfrak{m}_{B C}^{2} \circ \tilde{\lambda}_{C}\right)(g)(q)=\tilde{\lambda}_{C}(g)\left(m_{B C}^{2}(q)\right)= \\
& \tilde{\lambda}_{C}(g)\left(\sum_{i \in I} a_{i}\left(\sum_{k \in I} m_{i k}\right) z_{i}^{2}+\sum_{i<j \in I} b_{i j}\left[\left(\sum_{r \in I} m_{i r}\right)\left(\sum_{s \in I} m_{j s}\right)\right] z_{i} z_{j}\right)= \\
& \sum_{i \in I} a_{i}\left(\sum_{k \in I} m_{i k}\right) \alpha_{i}+\sum_{i<j \in I} b_{i j}\left[\left(\sum_{r \in I} m_{i r}\right)\left(\sum_{s \in I} m_{j s}\right)\right] \beta_{i j}
\end{aligned}
$$

On the other hand,

$$
\begin{aligned}
& \left(\tilde{\lambda}_{B} \circ \mu_{B C}^{2}\right)(g)(q)=\tilde{\lambda}\left(\mu_{B C}^{2}(g)(q)\right) \\
& \left.=\tilde{\lambda}_{B}\left(\left(t_{i}^{\alpha_{i}\left(\sum_{k \in I} m_{i k}\right)}\right)\left(t_{i j}^{\beta_{i j}} \sum_{r, s \in I} m_{i r} m_{j s}\right)\right)\left(1_{i}\right)\right)\left(\sum_{i \in I} a_{i} z_{i}^{2}+\sum_{i<j \in I} b_{i j} z_{i} z_{j}\right)= \\
& \sum_{i \in I} a_{i}\left(\alpha_{i}\left(\sum_{k \in I} m_{i k}\right)\right)+\sum_{i<j \in I} b_{i j}\left(\beta_{i j}\left(\sum_{r, s \in I} m_{i r} m_{j s}\right)\right)
\end{aligned}
$$

proving that $\mathfrak{m}_{B C}^{2} \circ \tilde{\lambda}_{C}=\tilde{\lambda}_{B} \circ \mu_{B C}^{2}$.
Lemma 7.2.9. With the notations of Lemmas 7.2.2 7.2 .7 we have the following.
$i$ - The arrows $\mu_{B C}^{1}$ and $\mu_{B C}^{2}$ are isomorphisms. Moreover, for all well-ordered basis $B, C, D$ we have $\mu_{B B}^{1}=i d, \mu_{B B}^{2}=i d, \mu_{B D}^{1}=\mu_{B C}^{1} \circ \mu_{C D}^{1}$ and $\mu_{B D}^{2}=\mu_{B C}^{2} \circ \mu_{C D}^{2}$.
$i i$ - The isomorphism $\mu_{B C}: \mathcal{W}(C) \rightarrow \mathcal{W}(B)$ restricts to an isomorphism $\mathcal{V}(C) \rightarrow \mathcal{V}(B)$ so we get quotient isomorphism

$$
\tilde{\mu}_{B C}: \mathcal{W}(C) / \mathcal{V}(C) \rightarrow \mathcal{W}(B) / \mathcal{V}(B) .
$$

iii - If $B, C, D$ are well-ordered base of $G$, then $\tilde{\mu}_{C C}=$ id and $\tilde{\mu}_{B D}=\tilde{\mu}_{B C} \circ \tilde{\mu}_{C D}$.
Proof.
i - Just use the same calculations made in Lemma 7.2.5,
ii - By Proposition 7.2.2 we have

$$
\mathfrak{m}_{B C}^{2}(Q(B))=Q(C) .
$$

Since $\mathcal{V}(B)=Q(B)^{\perp}$, we have an induced isomorphism $\left.\mu_{B C}\right|_{\mathcal{V}(B)}: \mathcal{V}(B) \xrightarrow{\cong} \mathcal{V}(C)$, legitimating the quotient isomorphism

$$
\tilde{\mu}_{B C}: \mathcal{W}(C) / \mathcal{V}(C) \rightarrow \mathcal{W}(B) / \mathcal{V}(B)
$$

iii - It is an immediate consequence of item (i).

Definition 7.2.10 (Galois Group - base independent version). Let $G$ be a pre-special group. Take

$$
E_{G}=\left\{B: B \text { is a well-ordered } \mathbb{F}_{2} \text {-basis of } G\right\} .
$$

Consider the set $E_{G}$ endowed with the trivial groupoid operation of concatenation of pairs (i.e., the arrows are $E_{G} \times E_{G}$ ) and take the functor Gal: $E_{G} \rightarrow \mathcal{C}$, (where $\mathcal{C}$ is the category of $\mathcal{C}$-groups) given by the following rules: for an object $B \in E_{G}, \operatorname{Gal}(B):=\mathcal{W}(B) / \mathcal{V}(B)$ and for an arrow $(B, C) \in E_{G}^{2}$,

$$
\operatorname{Gal}(B, C)=\mu_{B C}: \mathcal{W}(C) / \mathcal{V}(C) \rightarrow \mathcal{W}(B) / \mathcal{V}(B)
$$

We define the Galois group of $G$, notation $\operatorname{Gal}(G)$ by

$$
\operatorname{Gal}(G):=\lim _{B \in E_{G}} \mathcal{W}(B) / \mathcal{V}(B) .
$$

Remark 7.2.11. Keeping the notation above, note that

$$
\pi_{B}: \operatorname{Gal}(G):=\lim _{B \in E_{G}} \mathcal{W}(B) / \mathcal{V}(B) \rightarrow \mathcal{W}(B) / \mathcal{V}(B)
$$

is an isomorphism, for each $B \in E_{G}$. This holds because

$$
\operatorname{Gal}(B, C)=\mu_{B C}: \mathcal{W}(C) / \mathcal{V}(C) \rightarrow \mathcal{W}(B) / \mathcal{V}(B)
$$

is an isomorphism for each arrow $(B, C) \in E_{G}^{2}$.
It is desirable to achieve explicit calculations of $\operatorname{Gal}(G)$ for finite reduced special groups and boolean algebras.

### 7.3 On the structure of Galois Groups of Pre Special Groups

As occurs with fields, the Galois group of a pre-special groups (in a certain subclass) is able to encode many relevant quadratic information.

All pre-special groups occurring is this section will be assumed $k$-stable.
We start this Section, by providing more details on the structure of $\mathcal{C}$-groups.
Let $\mathcal{G}$ be a $\mathcal{C}$-group on $I$-minimal generators. By Theorem 7.1.11, There is an epimorphism $\lambda: \mathcal{W}(I) \rightarrow \mathcal{G}$ with kernel $\mathcal{V} \subseteq \Phi(I)$, and then, we have an isomorphism $\tilde{\lambda}: \mathcal{W}(I) / \mathcal{V} \rightarrow \mathcal{G}$.
fixhugo2
Proposition 7.3.1. With the above notation, we have the following.
$i$ - We have a natural bijection

$$
\begin{aligned}
& \{M \subseteq \mathcal{G}: M \text { is a maximal open subgroup }\} \cong \\
& \{M \subseteq \mathcal{W}(I): M \text { is a maximal open subgroup }\} .
\end{aligned}
$$

Then $\Phi(\mathcal{G}) \cong \Phi(I) / \mathcal{V}$ and $\mathcal{G} / \Phi(\mathcal{G}) \cong \mathcal{W}(I) / \Phi(I)$.
ii - The maximal closed subgroups of $\mathcal{W}(I)$ are precisely the clopen (normal) subgroups with quotient $\mathbb{Z}_{2}$. Moreover, we have a natural bijection

$$
\{M \subseteq \mathcal{G}: M \text { is a maximal open subgroup }\} \cong P_{\text {fin }}(I) \backslash\{\emptyset\} .
$$

and that extends to a natural bijection

$$
\{N \subseteq \mathcal{G}: N \text { is an open subgroup with index } \leq 2\} \cong P_{\text {fin }}(I)
$$

iii - We have a natural isomorphism of $\mathbb{Z}_{2}$-modules

$$
\operatorname{Homcont}\left(\mathcal{G}, \mathbb{Z}_{2}\right) \cong f \operatorname{sFunc}\left(I, \mathbb{Z}_{2}\right)
$$

where $f$ sFunc $\left(I, \mathbb{Z}_{2}\right)$ is the set of all function $f: I \rightarrow \mathbb{Z}_{2}$ with finite support.
Proof.
i - The desired bijection follows by the bijection

$$
\begin{aligned}
& \{M \subseteq \mathcal{W}(I) / \mathcal{V}: M \text { is a maximal open subgroup }\} \cong \\
& \{M \subseteq \mathcal{W}(I): M \text { is a maximal open subgroup }\} .
\end{aligned}
$$

given by the following rule: lets $q: \mathcal{W}(I) \rightarrow \mathcal{W}(I) / \mathcal{V}$ denote the canonical projection. We have a function $\bar{q}: \mathcal{P}(\mathcal{W}(I) / \mathcal{V}) \rightarrow \mathcal{P}(\mathcal{W}(I))$ given by the rule

$$
\bar{q}(X):=q^{-1}[X] \text { (the inverse image). }
$$

This function $\bar{q}$ induces the desired bijection. Since the Frattini subgroup of $\mathcal{G}$ is the intersection of all open normal subgroups we have (via $\tilde{\lambda}$ and the bijection) $\Phi(\mathcal{G}) \cong \Phi(I) / \mathcal{V}$. Then

$$
\mathcal{G} / \phi(\mathcal{G}) \cong(\mathcal{W}(I) / \mathcal{V}) /(\Phi(I) / \mathcal{V}) \cong \mathcal{W}(I) / \Phi(I) .
$$

ii - For $\left\{i_{0}, \ldots, i_{n}\right\} \subseteq I$ denote

$$
\zeta_{I}\left(i_{0}, \ldots, i_{n}\right):=\left\{\sigma \in \mathcal{W}(I): \gamma_{i_{0}}(\sigma)+\ldots+\gamma_{i_{n}}(\sigma)=0\right\} .
$$

We have that $\zeta_{I}\left(i_{0}, \ldots, i_{n}\right)$ is a subgroup of $\mathcal{W}(I)$. Now let $\tau, \theta \in \mathcal{W}(I) \backslash \zeta_{I}\left(i_{0}, \ldots, i_{n}\right)$. Then for all $i \in I$,

$$
\gamma_{i}\left(\theta^{-1} \tau\right)=\gamma_{i}(\theta)+\gamma_{i}(\tau)
$$

Therefore

$$
\sum_{p=1}^{n} \gamma_{i_{p}}\left(\theta^{-1} \tau\right)=\sum_{p=1}^{n}\left[\gamma_{i_{p}}(\theta)+\gamma_{i_{p}}(\tau)\right]=\sum_{p=1}^{n} \gamma_{i_{p}}(\theta)+\sum_{p=1}^{n} \gamma_{i_{p}}(\tau)=1+1=0 .
$$

Then $\theta^{-1} \tau \in \zeta_{I}\left(i_{0}, \ldots, i_{n}\right)$ which imply

$$
\mathcal{W}(I) / \zeta_{I}\left(i_{0}, \ldots, i_{n}\right)=\{\overline{1}, \bar{\tau}\} \cong \mathbb{Z}_{2}
$$

To verify that $\zeta_{I}\left(i_{0}, \ldots, i_{n}\right)$ is clopen, note that

$$
\gamma_{i_{0}}(\sigma)+\ldots+\gamma_{i_{n}}(\sigma)=0 \text { iff }\left\{\begin{array}{l}
\gamma_{i_{0}}(\sigma)+\ldots+\gamma_{i_{n-1}}(\sigma)=0 \text { and } \gamma_{i_{n}}(\sigma)=0 \text { or } \\
\gamma_{i_{0}}(\sigma)+\ldots+\gamma_{i_{n-1}}(\sigma)=1 \text { and } \gamma_{i_{n}}(\sigma)=1
\end{array}\right.
$$

In other words,

$$
\zeta_{I}\left(i_{0}, \ldots, i_{n}\right)=\left[\zeta_{I}\left(i_{0}, \ldots, i_{n-1}\right) \cap M_{i_{n}}\right] \cup\left[\zeta_{I}\left(i_{0}, \ldots, i_{n-1}\right)^{c} \cap M_{i_{n}}^{c}\right]
$$

So, in order to verify that $\zeta_{I}\left(i_{0}, \ldots, i_{n}\right)$ is clopen is enough to deal with the case $n=0$. But $\zeta_{I}\left(i_{0}\right)=M_{i_{0}}$ is in fact a clopen, which provide (by induction) that $\zeta_{I}\left(i_{0}, \ldots, i_{n}\right)$ is clopen for all $i_{0}, \ldots, i_{n} \in I$. Then $\zeta_{I}\left(i_{0}, \ldots, i_{n}\right)$ is a maximal clopen subgroup of $\mathcal{W}(I)$, and we have a well-defined injective function

$$
\zeta_{I}: \mathcal{P}_{\text {fin }}(I) \backslash\{\emptyset\} \rightarrow\{M \subseteq \mathcal{W}(I): M \text { is a maximal open subgroup }\}
$$

given by the rule $\left\{i_{0}, \ldots, i_{n}\right\} \mapsto \zeta_{I}\left(i_{0}, \ldots, i_{n}\right)$.
For surjectivity, let $M$ be a maximal open subgroup. Then $M$ is closed of finite index and by Lemma 7.1.9(a), $M$ has index 2 in $G$. Using Propositions 7.1 .5 and 7.1 .6 and the compacity of $\mathcal{W}(I)$, there exists $i_{1}, \ldots, i_{n}, j_{1}, \ldots, j_{m}, k_{1}, \ldots, k_{p} \in I$ with

$$
M_{i_{1}} \cap \ldots \cap M_{i_{n}} \cap S_{j_{1}} \cap \ldots \cap S_{j_{m}} \cap D_{k_{1}} \cap \ldots \cap D_{k_{p}} \subseteq M
$$

Note that we have

$$
M_{i_{1}} \cap \ldots \cap M_{i_{n}} \cap S_{j_{1}} \cap \ldots \cap S_{j_{m}} \cap D_{k_{1}} \cap \ldots \cap D_{k_{p}} \subseteq \zeta_{I}\left(i_{1}, \ldots, i_{n}, j_{1}, \ldots, j_{m}, k_{1}, \ldots, k_{p}\right)
$$

Lets denote $\zeta_{I}\left(i_{1}, \ldots, i_{n}, j_{1}, \ldots, j_{m}, k_{1}, \ldots, k_{p}\right):=\zeta_{I}(\vec{i}, \vec{j}, \vec{k})$ and $H:=M \cap \zeta_{I}(\vec{i}, \vec{j}, \vec{k})$. Suppose $H \neq M$ and let $\tau, \theta \in M \backslash H$. The same calculations made for injectivity shows that $\theta^{-1} \tau \in$ $\zeta_{I}(\vec{i}, \vec{j}, \vec{k})$ which imply

$$
M / H=\{\overline{1}, \bar{\tau}\} \cong \mathbb{Z}_{2}
$$

Moreover, using the same calculations made in Proposition 7.1.5(b) we have that $H$ has index $\mathbb{Z}_{4}$ in $\mathcal{W}(I)$. Since $H \subseteq M$ and $H \subseteq \zeta_{I}(\vec{i}, \vec{j}, \vec{k})$ with both maximal clopen subgroups, by Lemma 7.3.7(i) we have $M=\zeta_{I}(\vec{i}, \vec{j}, \vec{k})$, contradicting the assumption $H \neq M$. Then $H=M$ and we have $M=\zeta_{I}(\vec{i}, \vec{j}, \vec{k})$. Therefore we have bijections

$$
\begin{aligned}
\mathcal{P}_{f i n}(I) \backslash\{\emptyset\} & \cong\{M \subseteq \mathcal{W}(I): M \text { is a maximal open subgroup }\} \\
& \cong\{M \subseteq \mathcal{G}: M \text { is a maximal open subgroup }\}
\end{aligned}
$$

Therefore

$$
\{N \subseteq \mathcal{G}: N \text { is an open subgroup with index } \leq 2\} \cong P_{f i n}(I)
$$

iii - Since $\mathcal{G}=\mathcal{W}(I) / \mathcal{V}$ and

$$
\mathcal{V} \subseteq \Phi(I)=\bigcap\left\{k e r(\varphi): \varphi \in \operatorname{Homcont}\left(\mathcal{W}(I), \mathbb{Z}_{2}\right)\right\}
$$

then the natural epimorphism $\mathcal{W}(I) \rightarrow \mathcal{W}(I) / \mathcal{V}$ induces the isomorphism

$$
H \operatorname{cocont}\left(\mathcal{G}, \mathbb{Z}_{2}\right) \cong H \operatorname{comcont}\left(\mathcal{W}(I), \mathbb{Z}_{2}\right)
$$

Since $\mathbb{Z}_{2}$ is a finite/discrete $\mathcal{C}$-group, then the universal property of $\mathcal{W}(I)$ (Theorem 7.1.11) gives a natural isomorphism $\operatorname{Homcont}\left(\mathcal{W}(I), \mathbb{Z}_{2}\right) \cong f \operatorname{sFunc}\left(I, \mathbb{Z}_{2}\right)$.
Alternatively, the result follows also from the item (ii) above and the (obvious) natural bijections

$$
\begin{gathered}
\operatorname{Homcont}\left(\mathcal{G}, \mathbb{Z}_{2}\right) \cong\{N \subseteq \mathcal{G}: N \text { is an open subgroup with index } \leq 2\} \\
P_{\text {fin }}(I) \cong f \operatorname{sFunc}\left(I, \mathbb{Z}_{2}\right) .
\end{gathered}
$$

Proposition 7.3.2 (Prontryagin duality). Let $\mathcal{G}=\operatorname{Gal}(G)$ for some pre-special group $G$. Then
$i$ - There is a canonical bijection
$\mathbb{M}: G \cong\{M \subseteq \mathcal{G}: M$ is a (closed, normal) subgroup of index less or equal to 2$\}$
$a \mapsto M_{a}$ such that it induces a canonical bijection

$$
G \backslash\{1\} \cong\{M \subseteq \mathcal{G}: M \text { is a maximal subgroup }\} .
$$

ii - There is a canonical isomorphism of $\mathbb{Z}_{2}$-modules $\psi_{G}: G \xlongequal{\cong} \operatorname{Homcont}\left(\mathcal{G}, \mathbb{Z}_{2}\right), a \mapsto \mu_{a}$, where $\mu_{a}: \mathcal{G} \rightarrow \mathbb{Z}_{2}$ is the unique continuous homomorphism such that $\operatorname{ker}\left(\mu_{a}\right)=M_{a}$.
iii - There is a canonical isomorphism of pro-2-groups $\phi_{G}: \mathcal{G} / \Phi(\mathcal{G}) \rightarrow \operatorname{Hom}\left(G, \mathbb{Z}_{2}\right)$.
Proof.
i - This follows directly from Proposition 7.3.1. since $\pi_{B}: \operatorname{Gal}(G) \xlongequal{\cong} \mathcal{W}(B) / \mathcal{V}(B)$ and $\mathcal{V}(B) \subseteq$ $\Phi(B)$, for every well orderd basis $B$ of $G$.
ii - By Proposition 7.3.1 (iii), for each well ordered basis $B$ in $G$, we have a natural isomorphism of $\mathbb{Z}_{2}$-modules.

$$
\operatorname{Homcont}\left(\mathcal{W}(B) / \mathcal{V}(B), \mathbb{Z}_{2}\right) \cong \operatorname{fsFunc}\left(B, \mathbb{Z}_{2}\right)
$$

This is, in fact, an isomorphism of $\mathbb{Z}_{2}$-modules. Taking into account the isomorphisms of "change of base", we glue the above isomorphisms to obtain the natural isomorphism

$$
\operatorname{Homcont}\left(\operatorname{Gal}(G), \mathbb{Z}_{2}\right)=\operatorname{Homcont}\left(\lim _{B \in E_{G}} \mathcal{W}(B) / \mathcal{V}(B), \mathbb{Z}_{2}\right) \cong \lim _{B \in E_{G}} \operatorname{fSFunc}\left(B, \mathbb{Z}_{2}\right) \cong G
$$

iii - Note that $\pi_{\mathcal{G}}: \mathcal{G} \rightarrow \mathcal{G} / \Phi(\mathcal{G})$ induces a $\mathbb{Z}_{2}$-isomorphism

$$
\pi_{\mathcal{G}}^{*}: \operatorname{Homcont}\left(\mathcal{G} / \Phi(\mathcal{G}), \mathbb{Z}_{2}\right) \xlongequal{\cong} \operatorname{Homcont}\left(\mathcal{G}, \mathbb{Z}_{2}\right)
$$

By Lemma 7.2.7, the isomorphisms described in Proposition 7.1.15, namely

$$
\mathcal{W}(B) / \Phi(B) \cong \operatorname{Hom}\left(P_{1}(B), \mathbb{Z}_{2}\right),
$$

are natural. Thus we obtain a natural isomorphism

$$
\mathcal{G} / \Phi(\mathcal{G}) \cong \operatorname{Hom}\left(G, \mathbb{Z}_{2}\right)
$$

Remark 7.3.3. Note that combining items (iii), (ii) of the Proposition above, we obtain the Pontryagin duality:

$$
\begin{gathered}
\mathcal{G} / \Phi(\mathcal{G}) \cong \operatorname{Hom}\left(G, \mathbb{Z}_{2}\right) \cong \operatorname{Hom}\left(\operatorname{Homcont}\left(\mathcal{G} / \Phi(\mathcal{G}), \mathbb{Z}_{2}\right), \mathbb{Z}_{2}\right) \\
G \cong \operatorname{Hom} \operatorname{cont}\left(\mathcal{G} / \Phi(\mathcal{G}), \mathbb{Z}_{2}\right) \cong \operatorname{Homcont}\left(\operatorname{Hom}\left(G, \mathbb{Z}_{2}\right), \mathbb{Z}_{2}\right)
\end{gathered}
$$

This induces a canonical duality between the pointed $\mathbb{Z}_{2}$-module $(G,-1)$ and the "pointed" pro-2-group $\operatorname{Gal}(G), M)$, where $M \subseteq \operatorname{Gal}(G)$ is an open subgroup of index $\leq 2$.

Let $G$ be a k-stable special group. Write $\mathcal{G}=\operatorname{Gal}(G)$.
The isomorphism of pro-2-groups $\phi_{G}: \mathcal{G} / \Phi(\mathcal{G}) \rightarrow \operatorname{Hom}\left(G, \mathbb{Z}_{2}\right)$, determines a "perfect pairing" $\hat{\phi}_{G}: \mathcal{G} / \Phi(\mathcal{G}) \times G \rightarrow \mathbb{Z}_{2}$ given by the rule

$$
\hat{\phi}_{G}(\bar{\alpha}, g):=\langle\bar{\alpha}, g\rangle:=\phi_{G}(\bar{\alpha})(g) .
$$

We will denote ( $)^{\perp}$, generically, both the correspondences between subsets of $\mathcal{G} / \Phi(\mathcal{G})$ and subsets of $G$.
Proposition 7.3.4. The perfect pairing $\hat{\phi}_{G}: \mathcal{G} / \Phi(\mathcal{G}) \times G \rightarrow \mathbb{Z}_{2}$ gives an anti-isomorphism of complete lattices between the poset

$$
\begin{gathered}
\{R \subseteq \mathcal{G} / \Phi(\mathcal{G}): R \text { is a closed subgroup of } \mathcal{G} / \Phi(\mathcal{G})\}= \\
\{\pi(T) \subseteq \mathcal{G} / \Phi(\mathcal{G}): T \text { is a closed subgroup of } \mathcal{G}\}
\end{gathered}
$$

with $\pi: \mathcal{G} \rightarrow \mathcal{G} / \Phi(\mathcal{G})$ being the canonical projection, and the poset

$$
\{\Delta \subseteq G: \Delta \text { is a subgroup of } G\} .
$$

These anti-isomorphisms are given by the rules

$$
\begin{aligned}
\pi(T) & \mapsto \pi(T)^{\perp}:=\left\{a \in G: \hat{\phi}_{G}(\sigma / \Phi(G), a)=0 \text { for all } \sigma \in T\right\} \\
H & \mapsto H^{\perp}:=\left\{\sigma / \Phi(G): \hat{\phi}_{G}(\sigma / \Phi(G), a)=0 \text { for all } a \in H\right\} .
\end{aligned}
$$

From this we get:
$i$ - An anti-isomorphism of complete lattices between the posets

$$
\{\Delta \subseteq G: \Delta \text { is a subgroup of } G\}
$$

and

$$
\{T \subseteq \mathcal{G}: T \text { is a closed subgroup and } \Phi(\mathcal{G}) \subseteq T\}
$$

ii - A bijection between the sets

$$
\{\Delta \subseteq G: \Delta \text { is a maximal subgroup of } G\}
$$

and

$$
\{\pi(T) \subseteq \mathcal{G} / \Phi(\mathcal{G}): T \text { is a discrete subgroup with order } 2\}=
$$

$\{\pi(T) \subseteq \mathcal{G} / \Phi(\mathcal{G}): T$ is a closed subgroup with $\pi(T)=\{i d, \sigma / \Phi\}$, for some $\sigma \in \mathcal{G} \backslash \Phi(\mathcal{G})\}$.
Proof. All items are immediate consequences of the isomorphism.
Since $\mathcal{G}$ is a compact Hausdorff group and $\Phi(\mathcal{G}) \subseteq \mathcal{G}$ is a closed normal subgroup, note that then the quotient map $\mathcal{G} \rightarrow \mathcal{G} / \Phi(\mathcal{G})$ gives an isomorphism of complete lattices between the poset of closed subgroups of $\mathcal{G}$ which contains $\Phi(\mathcal{G})$ and the poset of closed subgroups of $\mathcal{G} / \Phi(\mathcal{G})$.

Remark 7.3.5. Let $\sigma \in \mathcal{G} \backslash \Phi(\mathcal{G})$. It follows from the definition of pairing that:

- For any $x \in G-\{1\}: \sigma \in M_{x}$ iff $<\sigma / \Phi(\mathcal{G}), x>=0$ iff $x \in\{\Phi(\mathcal{G}), \sigma . \Phi(\mathcal{G})\}^{\perp}$.
- If there is an involution in $\sigma . \Phi(\mathcal{G})$ then of all element in $\sigma . \Phi(\mathcal{G})$ are involutions.

To obtain quadratic information from the Galois groups, we will need develop deeper group theoretic results.

Zd-le
Lemma 7.3.6. Let $B$ an well ordered basis of $G$ and consider $\eta_{B}=\pi_{B}^{-1}: \mathcal{W}(B) / \mathcal{V}(B) \xlongequal{\cong} \operatorname{Gal}(G)$.
$i$ - Let $a \neq 1$, choose $B=\left\{a_{i}: i \in I\right\}$ an well ordered basis of $G$ such that $a \in B$, say $a=a_{i}$. Then $\eta_{B}\left[M_{i}^{\prime} / \mathcal{V}(B)\right]=M_{a}$.
ii- Let $a, b \neq 1, a \neq b$ so $\{a, b\}$ is $a \mathbb{Z}_{2}$-l.i. subset of $G$, choose $B=\left\{a_{i}: i \in I\right\}$ an well ordered basis of $G$ such that $a, b \in B$, say $a=a_{i}, b=a_{j}, i<j \in I$. Let $\left\{M_{i}^{\prime}, M_{j}^{\prime}, M^{\prime}\right\}$ be the three maximal subgroups of $\mathcal{W}(B)$ above $M_{i}^{\prime} \cap M_{j}^{\prime}$. Then $\eta_{B}\left[M_{i}^{\prime} / \mathcal{V}(B)\right]=M_{a}, \eta_{B}\left[M_{j}^{\prime} / \mathcal{V}(B)\right]=M_{b}$ and $\eta_{B}\left[M^{\prime} / \mathcal{V}(B)\right]=M_{a b}$.
iii- Let $\left\{M_{1}, M_{2}, M_{3}\right\} \subseteq \operatorname{Gal}(G)$ maximal subgroups that are pairwise distinct. Then are equivalent:

- $\left\{M_{1}, M_{2}, M_{3}\right\}$ are independent, which means that for each of 3 enumerations $\{u, v, w\}$ of $\{1,2,3\}, M_{u} \cap M_{v} \nsubseteq M_{w}$.
- There is some enumeration $\{u, v, w\}$ of $\{1,2,3\}$ with $M_{u} \cap M_{v} \nsubseteq M_{w}$.
- $\operatorname{Gal}(G) /\left(M_{1} \cap M_{2} \cap M_{3}\right) \cong \mathbb{Z}_{2} \times \mathbb{Z}_{2} \times \mathbb{Z}_{2}$

Proof.
i- Recall that $\mathcal{V}(B) \subseteq \Phi(B) \subseteq M_{k}^{\prime}, \forall k \in I, M_{i}^{\prime}=\left\{\sigma \in \mathcal{W}(B): \gamma_{i}(\sigma)=0\right\}$. The "isomorphic" perfect pairings,

$$
\begin{aligned}
& <,>_{B}: \mathcal{W}(B) / \Phi(B) \times P_{1}(B) \rightarrow \mathbb{Z}_{2} \\
& <,>: \operatorname{Gal}(G) / \Phi(\operatorname{Gal}(G)) \times G \rightarrow \mathbb{Z}_{2}
\end{aligned}
$$

provides $M_{i}^{\prime} / \Phi(B)=\left\{z_{i}\right\}^{\perp}$ and $M_{a_{i}} / \Phi(G a l(G))=\left\{a_{i}\right\}^{\perp}$. Since the pairings are "compatible", i.e., the dual of the isomorphism $P_{1}(B) \xlongequal{\cong} G: z_{k} \mapsto a_{k}, k \in I$ corresponds to the isomorphism

$$
\mathcal{W}(B) / \Phi(B) \underset{\text { can }}{\cong}(\mathcal{W}(B) / \mathcal{V}(B)) / \Phi((\mathcal{W}(B) / \mathcal{V}(B))) \underset{\vec{\eta}_{B}}{\cong} \operatorname{Gal}(G) / \Phi(\operatorname{Gal}(G)),
$$

we have $\bar{\eta}_{B} \circ \operatorname{can}\left[M_{i}^{\prime} / \Phi(B)\right]=M_{a_{i}} / \Phi(\operatorname{Gal}(G))$. Therefore, as the Frattini subgroups are contained in all maximal open subgroups, we get $\eta_{B}\left[M_{i}^{\prime} / \mathcal{V}(B)\right]=M_{a_{i}}$.
ii- $\quad M^{\prime}=\left\{\sigma \in \mathcal{W}(B): \gamma_{i}(\sigma)+\gamma_{j}(\sigma)=0\right\}$. The "isomorphic" perfect pairings,

$$
\mathcal{W}(B) / \Phi(B) \times P_{1}(B) \rightarrow \mathbb{Z}_{2} \text { and } \operatorname{Gal}(G) / \Phi(\operatorname{Gal}(G)) \times G \rightarrow \mathbb{Z}_{2}
$$

provides $M_{a} / \Phi(G a l(G))=\{a\}^{\perp}$ and $M_{b} / \Phi(G a l(G))=\{b\}^{\perp}$, so

$$
\begin{aligned}
M_{a b} / \Phi(\operatorname{Gal}(G)) & =\{a b\}^{\perp}=\{\theta / \Phi(\operatorname{Gal}(G)):<\theta / \Phi(\operatorname{Gal}(G)), a b>=0\} \\
& \left.\left.\subseteq\left(M_{a}\right) / \Phi(\operatorname{Gal}(G)) \cap M_{b}\right) / \Phi(\operatorname{Gal}(G))\right)=\left(M_{a} \cap M_{b}\right) / \Phi(\operatorname{Gal}(G))
\end{aligned}
$$

Therefore $\left\{M_{a}, M_{b}, M_{a b}\right\}$ are the three maximal subgroups of $\operatorname{Gal}(G)$ above $M_{a} \cap M_{b}$. Since $\left\{M_{i}^{\prime}, M_{j}^{\prime}, M^{\prime}\right\}$ are the three maximal subgroups of $\mathcal{W}(B)$ above $M_{i}^{\prime} \cap M_{j}^{\prime}$ and

$$
\eta_{B}\left[M_{i}^{\prime} / \mathcal{V}(B)\right]=M_{a} \text { and } \eta_{B}\left[M_{j}^{\prime} / \mathcal{V}(B)\right]=M_{b},
$$

we must have $\eta_{B}\left[M^{\prime} / \mathcal{V}(B)\right]=M_{a b}$
iii- We have uniquely determined $\{a, b, c\} \subseteq G \backslash\{1\}$, with $M_{1}=M_{a}, M_{2}=M_{b}, M_{3}=M_{c}$ and, from the hypothesis, $a, b, c$ are pairwise distinct so the result follows from (ii) since if $\{x, y\}$ is a $\mathbb{Z}_{2}$-li set then $\{x, x y\}$ and $\{y, x y\}$ are $\mathbb{Z}_{2}$-li sets and those 3 sets are $\mathbb{Z}_{2}$-basis of the group $\{1, x, y, x y\}$.

ZqDq-le
Lemma 7.3.7. Let $G$ be a $k$-stable pre-special group and denote $\mathcal{G}:=G a l(G)$.
$i$ - Let $S \subseteq \mathcal{G}$ a normal closed of with $\mathcal{G} / S \cong \mathbb{Z}_{4}$ then there is a unique maximal subgroup $H \subseteq \mathcal{G}$ such that $S \subseteq H$.
ii - Let $D \subseteq \mathcal{G}$ a normal closed of with $\mathcal{G} / S \cong \mathbb{D}_{4}$ then there is a unique set $\left\{H_{1}, H_{2}\right\}$ with $H_{1} \neq H_{2}, H_{i} \subseteq \mathcal{G}$ maximal subgroups such that $D \subseteq H_{1} \cap H_{2}$ and if $\left\{H, H_{1}, H_{2}\right\}$ are the three maximal subgroups above $H_{1} \cap H_{2}$ then $H / D \cong \mathbb{Z}_{4}$.

## Proof.

i - For the existential part take $r \in \mathcal{G} \backslash S$ such that $r^{2} \notin S$ and $\mathcal{G} / S=\left\{1 . S, r . S, r^{2} . S, r^{3} . S\right\} \cong \mathbb{Z}_{4}$ and take the maximal subgroup $H=1 S \cup r^{2} S$. Then $S \subseteq H$ and the canonical epimorphism $\mathcal{G} / S \rightarrow \mathcal{G} / H$ corresponds to the epimorphism $\mathbb{Z}_{4} \rightarrow \mathbb{Z}_{2}$. This maximal $H \supseteq S$ is unique because if $S \subseteq H_{1}, H_{2}$ with $H_{1} \neq H_{2}$, then $S \subseteq H_{1} \cap H_{2}$ and there will be an epimorphism $\mathbb{Z}_{4} \cong \mathcal{G} / S \rightarrow \mathcal{G} / H_{1} \cap H_{2} \cong \mathbb{Z}_{2} \times \mathbb{Z}_{2}$, however there is no element in $\mathbb{Z}_{2} \times \mathbb{Z}_{2}$ of order 4.
ii - For the existential part let $r, s \in \mathcal{G} \backslash D$ be such that $r^{2} \notin D, s^{2} \in D$ and then

$$
\mathbb{D}_{4} \cong \mathcal{G} / D=\left\{1 . D, r \cdot D, r^{2} . D, r^{3} . D, s \cdot D, s r \cdot D, s r^{2} \cdot D, s r^{3} . D\right\}
$$

Then each one of the maximal subgroups above $D$ is a reunion of four classes so they must contain $1 D, r^{2} D$. they are three:

$$
\begin{array}{r}
1 D \cup r^{2} D \cup r D \cup r^{3} D \\
1 D \cup r^{2} D \cup s D \cup s r^{2} D \\
1 D \cup r^{2} D \cup s r D \cup s r^{3} D
\end{array}
$$

Then take

$$
\begin{aligned}
\left\{H_{1}, H_{2}\right\} & =\left\{1 D \cup r^{2} D \cup s D \cup s r^{2} D, 1 D \cup r^{2} D \cup s r D \cup s r^{3} D\right\} \text { and } \\
H & =1 D \cup r^{2} D \cup r D \cup r^{3} D .
\end{aligned}
$$

Then $D \subseteq H_{1} \cap H_{2}=1 D \cup r^{2} D \subseteq H, H_{1}, H_{2}$. Then the canonical epimorphism $\mathcal{G} / D \rightarrow$ $\mathcal{G} / H_{1} \cap H_{2}$ corresponds to the epimorphism $\mathbb{D}_{4} \rightarrow \mathbb{Z}_{2} \times \mathbb{Z}_{2}$ and $H / D \cong \mathbb{Z}_{4}$. This pair of maximals $H_{1}, H_{2} \supseteq D$ is unique because if $H_{3}$ is a maximal $S \subseteq H_{3}$ with $H_{3} \neq H_{1}, H_{2}$, then $D \subseteq H_{1} \cap H_{2} \cap H_{3}$ and we have two to consider:

- $H_{1} \cap H_{2} \subseteq H_{3}$ : in this case $H_{3}=H=1 D \cup r^{2} D \cup r D \cup r^{3} D$, then

$$
1 D=\cup r^{2} D=H_{1} \cap H_{2}=H \cap H_{1}=H \cap H_{2}
$$

so $D \subseteq H \cap H_{1}$, but we see directly that $H_{2} / D \not \not \mathbb{Z}_{4}$; similarly $D \subseteq H \cap H_{2}$ but also $H_{1} / D \nexists \mathbb{Z}_{4}$;

- $H_{1} \cap H_{2} \nsubseteq H_{3}$ then $H_{3} \neq H, H_{1}, H_{2}$, also $H_{1} \cap H_{3} \nsubseteq H_{2}$ (because if $H_{1} \cap H_{3} \subseteq H_{2}$, then $H_{1} \cap H_{2} \nsubseteq H_{3}$, see the previous Lemma) and similarly $H_{2} \cap H_{3} \nsubseteq H_{1}$, then $\left\{H_{1}, H_{2}, H_{3}\right\}$ are independent, so in this case, the epimorphism $\mathcal{G} / D \rightarrow \mathcal{G} / H_{1} \cap H_{2} \cap H_{3}$ corresponds to an epimorphism $\mathbb{D}_{4} \rightarrow \mathbb{Z}_{2} \times \mathbb{Z}_{2} \times \mathbb{Z}_{2}$ but there is no element in $\mathbb{Z}_{2} \times \mathbb{Z}_{2} \times \mathbb{Z}_{2}$ of order 4 .


## Theorem 7.3.8.

$i$ - Let $a \in G \backslash\{1\}$.

$$
l(a) . l(a)=0 \in k_{2}(G) \Rightarrow
$$

There is $S \subseteq \operatorname{Gal}(G)$, a normal clopen subgroup such that $\operatorname{Gal}(G) / S \cong \mathbb{Z}_{4}$ and $S \subseteq M_{a}$.
ii - Let $a, b \in G \backslash\{1\}$ such that $a \neq b$

$$
l(a) \cdot l(b)=0 \in k_{2}(G) \Rightarrow
$$

There is $D \subseteq \operatorname{Gal}(G)$, a normal clopen subgroup such that $\operatorname{Gal}(G) / D \cong \mathbb{D}_{4}$ and $D \subseteq M_{a} \cap M_{b}$, $M_{a b} / D \cong \mathbb{Z}_{4}$.

Proof.
i - Since $\{a\} \subseteq G$ is a l.i. subset, take an well ordered basis $B=\left\{a_{i}: i \in I\right\} \subseteq G$ such that $a \in B$, say $a=a_{i}$. Then $\eta_{B}: \mathcal{W}(B) / \mathcal{V}(B) \xlongequal{\cong} \operatorname{Gal}(G)$ with $\mathcal{V}(B) \subseteq \Phi(B)$ and $\mathcal{V}(B)=Q(B)^{\perp}$ where $Q(B)=k e r\left(P_{2}(B) \rightarrow k_{2}(G)\right)$, then, by Proposition 7.2.2. $Q(B)=\left[\left\{q_{x y}^{B}: l(x) l(y)=0\right\}\right]$ so

$$
\mathcal{V}(B)=\bigcap\left\{\left(q_{x y}^{B}\right)^{\perp}: l(x) l(y)=0\right\}=\bigcap\left\{\left(q_{x y}^{B}\right)^{\perp}: x, y \neq 1, l(x) l(y)=0\right\} .
$$

Denote $M_{i}^{\prime}=\left\{\sigma \in \mathcal{W}(B): \gamma_{i}(\sigma)=0\right\}$ and $S_{i}^{\prime}=\left\{\sigma \in \mathcal{W}(B): \alpha_{i}(\sigma)=\gamma_{i}(\sigma)=0\right\}$ then, by Proposition 7.1.5(i), $S_{i}^{\prime} \subseteq M_{i}^{\prime} \subseteq \mathcal{W}(B)$ are clopen normal subgroups with $\mathcal{W}(B) / M_{i}^{\prime} \cong \mathbb{Z}_{2}$, $\mathcal{W}(B) / S_{i}^{\prime} \cong \mathbb{Z}_{4}$. We have $\mathcal{V}(B) \subseteq \Phi(B) \subseteq M_{i}^{\prime}$, and we state the

Claim: $\mathcal{V}(B) \subseteq S_{i}^{\prime}$.
This entails that $M_{a}=\eta_{B}\left[M_{i}^{\prime} / \mathcal{V}(B)\right] \subseteq \operatorname{Gal}(G), \operatorname{Gal}(G) / M_{a} \cong \mathbb{Z}_{2}$ and

$$
S:=\eta_{B}\left[S_{i}^{\prime} / \mathcal{V}(B)\right] \subseteq \operatorname{Gal}(G)
$$

is a clopen normal subgroup of $\operatorname{Gal}(G)$ with $\operatorname{Gal}(G) / S \cong \mathbb{Z}_{4}$ and $S \subseteq M_{a}$, as we need.

Proof of the Claim: We will see that $S_{i}^{\prime} \cap \Phi(B)=\left(q_{a_{i} a_{i}}^{B}\right)^{\perp}$ then as $a=a_{i}$ and $l(a) l(a)=0$ we get $\mathcal{V}(B) \subseteq\left(q_{a_{i} a_{i}}^{B}\right)^{\perp}$ so $\mathcal{V}(B) \subseteq S_{i}^{\prime} \cap \Phi(B)$. Since $1 \neq a=a_{i}$ it follows that $q_{a_{i} a_{i}}^{B}=z_{i}^{2} \in P_{2}(B)$ is such that $\left(q_{a_{i} a_{i}}^{B}\right)^{\perp} \subseteq \Phi(B)$ has index 2 and we will proof that $S_{i}^{\prime} \cap \Phi(B) \subseteq \Phi(B)$ has also index 2 and $S_{i}^{\prime} \cap \Phi(B) \subseteq\left(q_{a_{i} a_{i}}^{B}\right)^{\perp}$ so we get $S_{i}^{\prime} \cap \Phi(B)=\left(q_{a_{i} a_{i}}^{B}\right)^{\perp}$. Firstly we show that $\Phi(B) / S_{i}^{\prime} \cap \Phi(B) \cong \mathbb{Z}_{2}$ : as $\Phi(B) \hookrightarrow M_{i}^{\prime}$ then $\Phi(B) / S_{i}^{\prime} \cap \Phi(B) \mapsto M_{i}^{\prime} / S_{i}^{\prime}$ and $M_{i}^{\prime} / S_{i}^{\prime} \cong \mathbb{Z}_{2}$ so $\Phi(B) / S_{i}^{\prime} \cap \Phi(B)$ has 1 or 2 elements. However, it cannot has 1 element: if $S_{i}^{\prime} \cap \Phi(B)=\Phi(B)$ then, by Proposition 7.1.13,

$$
\bigcap\left\{M_{j}^{\prime}: j \in I\right\}=\Phi(B) \subseteq S_{i}^{\prime}
$$

but $\mathcal{W}(B)$ is a compact space and $S_{i}^{\prime} \subseteq \mathcal{W}(B)$ is open subset, $M_{j}^{\prime} \subseteq \mathcal{W}(B)$ is a closed subset $j \in I$ so there is a finite subset $\left\{j_{0}, \ldots, j_{n}\right\} \subseteq I$ such that $M_{j_{0}}^{\prime} \cap \ldots \cap M_{j_{n}}^{\prime} \subseteq S_{i}^{\prime}$, choose $n \in \mathbb{N}$ minimum with this property so for each $m \leq n, \bigcap\left\{M_{j_{l}}^{\prime}: l \neq m\right\} \nsubseteq M_{j_{m}}^{\prime}$ then we have an isomorphism

$$
\operatorname{Gal}(G) / \bigcap\left\{M_{j_{l}}^{\prime}: l \leq n\right\} \xlongequal{\cong} \prod_{l \leq n} \operatorname{Gal}(G) / M_{j_{l}}^{\prime}
$$

so the epimorphim $\operatorname{Gal}(G) / \bigcap\left\{M_{j_{l}}^{\prime}: l \leq n\right\} \rightarrow \operatorname{Gal}(G) / S_{i}^{\prime}$ corresponds to an epimorphism $\prod_{l \leq n} \mathbb{Z}_{2} \rightarrow \mathbb{Z}_{4}$, but the two elements of order 4 in $\mathbb{Z}_{4}$ cannot be in the image of the homomorphism. Now we prove that $S_{i}^{\prime} \cap \Phi(B) \subseteq\left(q_{a_{i} a_{i}}^{B}\right)^{\perp}$ : we have $\left(q_{a_{i} a_{i}}^{B}\right)^{\perp}=\left\{z_{i}^{2}\right\}^{\perp}$ and

$$
S_{i}^{\prime} \cap \Phi(B)=\left\{\sigma \in \mathcal{W}(I): \alpha_{i}(\sigma)=0 \text { and } \gamma_{j}(\sigma)=0 \text { for each } j \in I\right\}
$$

and it follows from of the group operation and the definition of the pairing $<,>: \Phi(B) \times$ $P_{2}(B) \rightarrow \mathbb{Z}_{2}$ that

$$
\left\{x_{k} x_{l}: k<l \in I\right\} \cup\left\{x_{j}^{2}: i \neq j \in I\right\} \subseteq\left(S_{i}^{\prime} \cap \Phi(B)\right) \cap\left\{z_{i}^{2}\right\}^{\perp}
$$

Since $S_{i}^{\prime} \cap \Phi(B),\left\{z_{i}^{2}\right\}^{\perp} \subseteq \Phi(B)$ are closed subgroups,

$$
\text { closure }\left(\left[\left\{x_{k} x_{l}: k<l \in I\right\} \cup\left\{x_{j}^{2}: i \neq j \in I\right\}\right]\right) \subseteq\left(S_{i}^{\prime} \cap \Phi(B)\right) \cap\left\{z_{i}^{2}\right\}^{\perp} .
$$

Now we will prove that $S_{i}^{\prime} \cap \Phi(B) \subseteq \operatorname{closure}\left(\left[\left\{x_{k} x_{l}: k<l \in I\right\} \cup\left\{x_{j}^{2}: i \neq j \in I\right\}\right]\right)$; it is enough find for each $\sigma \in S_{i}^{\prime} \cap \Phi(B)$ and each basic neighborhood $T$ of $1 \in \mathcal{W}(B)$ two finite sets $\left\{j_{1}, \ldots j_{n}\right\} \subseteq I-\{i\}$ and $\left\{\left(k_{1}, l_{1}\right), \ldots\left(k_{m}, l_{m}\right): k_{u}<l_{u} \in I, 1 \leq u \leq m\right\}$ such that $\left(x_{j_{1}}^{2} \ldots \ldots x_{j_{n}}^{2} \cdot x_{k_{1}} x_{l_{1}} \ldots . x_{k_{m}} \cdot x_{l_{m}}\right) \in \sigma \cdot T$ : let $T=\bigcap U$ where

$$
U \subseteq_{f i n} V=\left\{M_{j}^{\prime}: j \in I\right\} \cup\left\{S_{j}^{\prime}: j \in I\right\} \cup\left\{D_{k l}^{\prime}: k<l \in I\right\}
$$

and take

$$
\begin{aligned}
\left\{j_{1}, \ldots j_{n}\right\} & =\left\{j \in I: S_{j}^{\prime} \in U, \alpha_{j}(\sigma)=1\right\} \subseteq I \backslash\{i\} \text { and } \\
\left\{k_{1}<l_{1}, \ldots k_{m}<l_{m}\right\} & =\left\{k<l \in I: D_{k l}^{\prime} \in U, \beta_{k l}(\sigma)=1\right\}
\end{aligned}
$$

$\left\{j_{1}, \ldots j_{n}\right\}=\left\{j \in I: S_{j}^{\prime} \in U, \alpha_{j}(\sigma)=1\right\} \subseteq I \backslash\{i\}$ and $\left\{k_{1}<l_{1}, \ldots k_{m}<l_{m}\right\}=\{k<l \in I:$ $\left.D_{k l}^{\prime} \in U, \beta_{k l}(\sigma)=1\right\}$ then, since $\gamma_{j}(\sigma)=0$, for all $j \in I$, we get

$$
\left(x_{j_{1}}^{2} \ldots \ldots x_{j_{n}}^{2} \cdot x_{k_{1}} x_{l_{1}} \ldots \ldots x_{k_{m}} \cdot x_{l_{m}}\right) \in \sigma \cdot T
$$

ii - Since $\{a, b\} \subseteq G$ is a l.i. subset, take an well ordered basis $B=\left\{a_{i}: i \in I\right\} \subseteq G$ such that $a, b \in B$, say $a=a_{i}, b=a_{j}, i<j \in I$. We denote

$$
\begin{aligned}
M_{i}^{\prime} & =\left\{\sigma \in \mathcal{W}(B): \gamma_{i}(\sigma)=0\right\} \\
M_{j}^{\prime} & =\left\{\sigma \in \mathcal{W}(B): \gamma_{j}(\sigma)=0\right\} \\
M^{\prime} & =\left\{\sigma \in \mathcal{W}(B): \gamma_{i}(\sigma)+\gamma_{j}(\sigma)=0\right\} \\
D_{i j}^{\prime} & =\left\{\sigma \in \mathcal{W}(B): \beta_{i j}(\sigma)=\gamma_{i}(\sigma)=\gamma_{j}(\sigma)=0\right\}
\end{aligned}
$$

By Proposition 7.1 .5 (ii), $D_{i j}^{\prime} \subseteq M_{i}^{\prime}, M_{j}^{\prime} \subseteq \mathcal{W}(B)$ are clopen normal subgroups with $\mathcal{W}(B) / M_{i}^{\prime} \cap$ $M_{j}^{\prime} \cong \mathbb{Z}_{2} \times \mathbb{Z}_{2}, \mathcal{W}(B) / D_{i j}^{\prime} \cong \mathbb{D}_{4}$. Besides $\mathcal{V}(B) \subseteq \Phi(B) \subseteq M_{i}^{\prime} \cap M_{j}^{\prime}$ and we stat the

Claim: $\mathcal{V}(B) \subseteq D_{i j}^{\prime}$.
This entails that $M_{a}=\eta_{B}\left[M_{i}^{\prime} / \mathcal{V}(B)\right] \subseteq G a l(G), M_{b}=\eta_{B}\left[M_{j}^{\prime} / \mathcal{V}(B)\right], M_{a b}=\eta_{B}\left[M^{\prime} / \mathcal{V}(B)\right]$, $\operatorname{Gal}(G) / M_{a} \cong \mathbb{Z}_{2} \cong \operatorname{Gal}(G) / M_{b}$ and $D \doteq \eta_{B}\left[S_{i}^{\prime} / \mathcal{V}(B)\right] \subseteq G a l(G)$ is a clopen normal subgroup of $G a l(G)$ with $G a l(G) / D \cong \mathbb{D}_{4}$ such that $D \subseteq M_{a} \cap M_{b}$ and $M_{a b} / D \cong \mathbb{Z}_{4}$, as we need.

Proof of the Claim: We will see that $D_{i j}^{\prime} \cap \Phi(B)=\left(q_{a_{i} a_{j}}^{B}\right)^{\perp}$ then as $a=a_{i}, b=a_{j}$ and $l(a) l(b)=0$ we get $\mathcal{V}(B) \subseteq\left(q_{a_{i} a_{j}}^{B}\right)^{\perp}$ so $\mathcal{V}(B) \subseteq D_{i j}^{\prime} \cap \Phi(B)$. As $1 \neq a=a_{i}$ and $1 \neq b=a_{j}$ with $i<j \in I$, it follows that $q_{a_{i} a_{j}}^{B}=z_{i} z_{j} \in P_{2}(B)$ is such that $\left(q_{a_{i} a_{j}}^{B}\right)^{\perp} \subseteq \Phi(B)$ has index 2 and we will proof that $D_{i j}^{\prime} \cap \Phi(B) \subseteq \Phi(B)$ has also index 2 and $D_{i j}^{\prime} \cap \Phi(B) \subseteq\left(q_{a_{i} a_{j}}^{B}\right)^{\perp}$ so we get $D_{i j}^{\prime} \cap \Phi(B)=\left(q_{a_{i} a_{j}}^{B}\right)^{\perp}$.

Firstly, we will prove that $\Phi(B) / D_{i j}^{\prime} \cap \Phi(B) \cong \mathbb{Z}_{2}$ : since $\Phi(B) \hookrightarrow M_{i}^{\prime} \cap M_{j}^{\prime}$ then

$$
\Phi(B) / D_{i j}^{\prime} \cap \Phi(B) \longmapsto M_{i}^{\prime} \cap M_{j}^{\prime} / D_{i j}^{\prime} \text { and } M_{i}^{\prime} \cap M_{j}^{\prime} / D_{i j}^{\prime} \cong \mathbb{Z}_{2}
$$

so $\Phi(B) / D_{i j}^{\prime} \cap \Phi(B)$ has 1 or 2 elements. However it cannot has 1 element: if $D_{i j}^{\prime} \cap \Phi(B)=$ $\Phi(B)$ then $\bigcap\left\{M_{k}^{\prime}: k \in I\right\}=\Phi(B) \subseteq D_{i j}^{\prime}$ but $\mathcal{W}(B)$ is a compact space and $D_{i j}^{\prime} \subseteq \mathcal{W}(B)$ is open subset,$M_{k}^{\prime} \subseteq \mathcal{W}(B)$ is a closed subset $k \in I$ so there is a finite subset $\left\{k_{0}, \ldots, k_{n}\right\} \subseteq I$ such that $M_{k_{0}}^{\prime} \cap \ldots \cap M_{k_{n}}^{\prime} \subseteq D_{i j}^{\prime}$, choose $n \in \mathbb{N}$ minimum with this property so for each $m \leq n$ , $\bigcap\left\{M_{k_{l}}^{\prime}: l \neq m\right\} \nsubseteq M_{k_{m}}^{\prime}$ then we have an isomorphism

$$
G a l(G) / \bigcap\left\{M_{k_{l}}^{\prime}: l \leq n\right\} \stackrel{\Im}{\rightrightarrows} \prod_{l \leq n} G a l(G) / M_{k_{l}}^{\prime}
$$

so the epimorphism $\operatorname{Gal}(G) / \bigcap\left\{M_{k_{l}}^{\prime}: l \leq n\right\} \rightarrow \operatorname{Gal}(G) / D_{i j}^{\prime}$ corresponds to an epimorphism $\prod_{l \leq n} \mathbb{Z}_{2} \rightarrow \mathbb{D}_{4}$, but the two elements of order 4 in $\mathbb{D}_{4}$ cannot be in the image of the homomorphism.

Now we prove that $D_{i j}^{\prime} \cap \Phi(B) \subseteq\left(q_{a_{i} a_{j}}^{B}\right)^{\perp}$ : we have $\left(q_{a_{i} a_{j}}^{B}\right)^{\perp}=\left\{z_{i} z_{j}\right\}^{\perp}$ and

$$
D_{i j}^{\prime} \cap \Phi(B)=\left\{\sigma \in \mathcal{W}(I): \beta_{i j}(\sigma)=0 \text { and } \gamma_{k}(\sigma)=0 \text { for each } k \in I\right\}
$$

and it follows from of the group operation and the definition of the pairing

$$
<,>: \Phi(B) \times P_{2}(B) \rightarrow \mathbb{Z}_{2}
$$

that

$$
\left\{x_{k} x_{l}: k<l \in I,(k, l) \neq(i, j)\right\} \cup\left\{x_{k}^{2}: k \in I\right\} \subseteq\left(D_{i j}^{\prime} \cap \Phi(B)\right) \cap\left\{z_{i} z_{j}\right\}^{\perp}
$$

then, since $D_{i j}^{\prime} \cap \Phi(B),\left\{z_{i} z_{j}\right\}^{\perp} \subseteq \Phi(B)$ are closed subgroups,

$$
\operatorname{closure}\left(\left[\left\{x_{k} x_{l}: k<l \in I,(k, l) \neq(i, j)\right\} \cup\left\{x_{k}^{2}: k \in I\right\}\right]\right) \subseteq\left(D_{i j}^{\prime} \cap \Phi(B)\right) \cap\left\{z_{i} z_{j}\right\}^{\perp}
$$

Now we will prove that $D_{i j}^{\prime} \cap \Phi(B) \subseteq \operatorname{closure}\left(\left[\left\{x_{k} x_{l}: k<l \in I,(k, l) \neq(i, j)\right\} \cup\left\{x_{k}^{2}: k \in I\right\}\right]\right)$; it is enough find for each $\sigma \in D_{i j}^{\prime} \cap \Phi(B)$ and each basic neighborhood $T$ of $1 \in \mathcal{W}(B)$ two finite sets $\left\{j_{1}, \ldots, j_{n}\right\} \subseteq I$ and $\left\{\left(k_{1}, l_{1}\right), \ldots\left(k_{m}, l_{m}\right): k_{u}<l_{u} \in I,\left(k_{u}, l_{u}\right) \neq(i, j), 1 \leq u \leq m\right\}$ such that $\left(x_{j_{1}}^{2} \ldots x_{j_{n}}^{2} \cdot x_{k_{1}} x_{l_{1}} \ldots . x_{k_{m}} \cdot x_{l_{m}}\right) \in \sigma . T:$ let $T=\bigcap U$ where

$$
U \subseteq_{f i n} V=\left\{M_{j}^{\prime}: j \in I\right\} \cup\left\{S_{j}^{\prime}: j \in I\right\} \cup\left\{D_{k l}^{\prime}: k<l \in I\right\}
$$

and take

$$
\begin{aligned}
\left\{j_{1}, \ldots j_{n}\right\} & =\left\{j \in I: S_{j}^{\prime} \in U, \alpha_{j}(\sigma)=1\right\} \subseteq I \text { and } \\
\left\{k_{1}<l_{1}, \ldots k_{m}<l_{m}\right\} & =\left\{k<l \in I: D_{k l}^{\prime} \in U, \beta_{k l}(\sigma)=1\right\} \subseteq I \times I \backslash\{(i, j)\}
\end{aligned}
$$

then, since $\gamma_{k}(\sigma)=0$, for all $k \in I$, we get

$$
\left(x_{j_{1}}^{2} \ldots x_{j_{n}}^{2} \cdot x_{k_{1}} x_{l_{1}} \ldots x_{k_{m}} \cdot x_{l_{m}}\right) \in \sigma . T
$$

The above proposition suggests the following:
pSGstandard-def
Definition 7.3.9. A pre-special group $G$ is said to be standard if it is a $k$-stable pre-special group and holds both the reverse implications in the Theorem 7.3.8 above.
Remark 7.3.10. Lemma 7.3.7 determines (injective) maps

$$
\begin{aligned}
& j_{1}:\left\{S \subseteq \mathcal{G}: S \text { is a normal subgroup of index } \mathbb{Z}_{4}\right\} \rightarrow \\
& \quad\{M \subseteq \mathcal{G}: M \text { is a maximal subgroup }\} \\
& j_{2}:\left\{D \subseteq \mathcal{G}: D \text { is a normal subgroup of index } \mathbb{D}_{4}\right\} \rightarrow \\
& \quad\left\{\left\{M_{1}, M_{2}\right\}: M_{1}, M_{2} \subseteq \mathcal{G}, M_{1} \neq M_{2} \text { are maximal subgroups }\right\}
\end{aligned}
$$

By the canonical bijection $\mathbb{M}: G \backslash\{1\} \stackrel{\cong}{\cong}\{M \subseteq G a l(G): M$ is a maximal clopen subgroup $\}$, it is natural to ask:
(1) Which subset of $\{a \in G: a \neq 1\}$ corresponds bijectively with image $\left(j_{1}\right)$ ?
(2) Which subset of $\{\{a, b\} \subseteq G: a, b \neq 1, a \neq b\}$ corresponds bijectively with image $\left(j_{2}\right)$ ?

The concept of standard pre-special group provides a full answer to these questions:
(1) The set $\left\{\{a\} \subseteq G:\{a\} l . i ., l(a) l(a)=0 \in k_{2}(G)\right\}$ corresponds bijectively with image $\left(j_{1}\right)$.
(2) The set $\left\{\{a, b\} \subseteq G:\{a, b\} l . i ., l(a) l(b)=0 \in k_{2}(G)\right\}$ corresponds bijectively with image $\left(j_{2}\right)$.

It follows from Propositions 2.3 and 2.4 in [53] that $S G(F)$ is a standard special group, for every field $F$ with $\operatorname{char}(F) \neq 2$.

We have already established that every special group $G$ is $k$-stable (see Proposition 6.3.7(iii)).
These suggest the following:

Question 7.3.11. Is every special group $G$ standard 1 田
In the sequel, we will see the relevance of the subclass of standard pre-special groups. We invite the reader to recall Proposition 7.3.4.

Theorem 7.3.12. Let $G$ be a $k$-stable pre-special group and denote $\mathcal{G}:=\operatorname{Gal}(G)$.
$i$ - Let $\sigma \in \mathcal{G} \backslash \Phi(\mathcal{G})$ be such that $\sigma^{2}=i d$. Then $\{\Phi(\mathcal{G}), \sigma . \Phi(\mathcal{G})\}^{\perp} \subseteq G$ is a maximal saturated subgroup of $G$.
ii- Suppose that $G$ is a standard pre-special group. Let $\sigma \in \mathcal{G} \backslash \Phi(\mathcal{G})$ be such that $\sigma^{2} \neq i d$ (so $\left.\sigma^{4}=i d\right)$. Then $\{\Phi(\mathcal{G}), \sigma . \Phi(\mathcal{G})\}^{\perp} \subseteq G$ is not a saturated subgroup of $G$.
iii- Suppose that $G$ is a standard special group. The set of all classes of conjugacy of involutions $\sigma \in \mathcal{G} \backslash \Phi(G)$ corresponds to the set of all orderings (= maximal saturated subgroups) of $G$.
$i v$ - Suppose that $G$ is a standard pre-special group. Then we have an anti-isomorphism of complete lattices between the posets
$\{\Delta \subseteq G: \Delta$ is a saturated subgroups of $G\}$
and
$\{T \subseteq \mathcal{G}: T$ is a closed subgroup of $\mathcal{G}($ topologically) generated by involutions such that $\Phi(\mathcal{G}) \subseteq T\}$

## Proof.

i- Let $\bar{T}=\{\Phi(\mathcal{G}), \sigma \Phi(\mathcal{G})\}$.

Claim: To have that $\bar{T}^{\perp} \subseteq G$ is a saturated subgroup it is enough to prove the following: $\forall x, y \in G$ if $\langle x, y\rangle \equiv<1, x y\rangle$, then $x \in \bar{T}^{\perp}$ or $y \in \bar{T}^{\perp}$.

Proof of Claim: Firstly we prove that $-1 \notin \bar{T}^{\perp}$ : take any $x \notin \bar{T}^{\perp}$ (there is some $x$, as $\bar{T}^{\perp} \subseteq G$ has index 2) then as $\langle x,-x>\equiv<1,-1>$ it follows from assumption in the claim that $-x \in \bar{T}^{\perp}$ so if $-1 \in \bar{T}^{\perp}$ then $x=-1 .(-x) \in \bar{T}^{\perp}$, a contradiction. Now let us prove that $\bar{T}^{\perp}$ is saturated: take any $a, b \in G$ such that $\left.b \in D_{G}(<1, a\rangle\right)$, assume $a \in \bar{T}^{\perp}$ then we have to prove that $b \in \bar{T}^{\perp}$ : as $(-a) \cdot a=-1 \notin \bar{T}^{\perp}$ then $-a \notin \bar{T}^{\perp}$ and as $\langle b, b a\rangle \equiv\langle 1, a\rangle$ we have $\langle b,-a\rangle \equiv<1,-b a\rangle$ so, by the assumption in the claim, we get $b \in \bar{T}^{\perp}$.
Now we will prove that $\sigma^{2}=i d$ entails $\forall x, y \in G$ if $\langle x, y\rangle \equiv<1, x y>$ then $x \in \bar{T}^{\perp}$ or $y \in \bar{T}^{\perp}$ :
We have three cases:

* $x$ (or $y$ ) is 1 ;
* $x=y \neq 1$;
* $x, y \neq 1$ and $x \neq y$.

There is nothing to proof in the first case. Now consider $x \in G \backslash\{1\}$ such that $\langle x, x\rangle \equiv<$ $1,1>$ : we must prove that $x \in \bar{T}^{\perp}$. Since $G$ is $k$-stable, we have $l(x) l(x)=l(1) l(1)=0$ then, by Theorem 7.3.8(i), there is a $S \subseteq \mathcal{G}$ a clopen normal subgroup such that $\mathcal{G} / S \cong \mathbb{Z}_{4}$ and

[^16]$S \subseteq M_{x}$. Consider the quotient map $p_{S}: \mathcal{G} \rightarrow \mathcal{G} / S$ and write $\mathcal{G} / S=\left\{1 / S, r / S, r^{2} / S, r^{3} / S\right\}$ then as $\sigma^{2}=i d$ we must have $\sigma / S \in\left\{1 / S, r^{2} / S\right\}$. If $\sigma / S=1 / S$ then $\sigma \in S \subseteq M_{x}$ i.e. $<\sigma / \Phi(\mathcal{G}), x>=0$ so $x \in\{\Phi(\mathcal{G}), \sigma \Phi(\mathcal{G})\}^{\perp}$. If $\sigma / S=r^{2} / S$ then $\sigma . r^{2}=\sigma . r^{-2} \in S \subseteq M_{x}$ i.e. $<\left(\sigma \cdot r^{2}\right) / \Phi(\mathcal{G}), x>=0$ but
$$
<\left(\sigma \cdot r^{2}\right) / \Phi(\mathcal{G}), x>=<\sigma / \Phi(\mathcal{G}), x>+<r / \Phi(\mathcal{G}), x>+<r / \Phi(\mathcal{G}), x>=<\sigma / \Phi(\mathcal{G}), x>
$$
then $<\sigma / \Phi(\mathcal{G}), x>=0$ so $x \in\{\Phi(\mathcal{G}), \sigma \Phi(\mathcal{G})\}^{\perp}$. Now take $x, y \in G \backslash\{1\}$ with $x \neq y$ and $<x, y>\equiv<1, x y>$ and supose $x \notin \bar{T}^{\perp}$ then we must prove that $y \in \bar{T}^{\perp}$. As $\{x, y\}$ is a two element l.i. set and $l(x) l(y)=l(1) l(x y)=0$ we have, by Theorem 7.3.8(ii), some $D \subseteq \mathcal{G}$ a clopen normal subgroup such that $\mathcal{G} / D \cong \mathbb{D}_{4}, D \subseteq M_{x} \cap M_{y}$ and $M_{x . y} / D \cong \mathbb{Z}_{4}$. Consider the quotient morphism $p_{D}: \mathcal{G} \rightarrow \mathcal{G} / D$ and write
$$
\mathcal{G} / D=\left\{1 / D, r / D, r^{2} / D, r^{3} / D, s / D, s r / D, s r^{2} / D, s r^{3} / D\right\}
$$
then, as $\sigma^{2}=i d$, we have $\sigma / D \notin\left\{r / D, r^{3} / D\right\}$. Let us prove that $\sigma / D \notin\left\{1 / D, r^{2} / D\right\}$ : as $<1 / \Phi(\mathcal{G}), x>=<r^{2} / \Phi(\mathcal{G}), x>=0$ we have $\left\{i d, r^{2}\right\} \subseteq M_{x}$ and as we selected $x \notin\{\Phi, \sigma \Phi\}^{\perp}$ we have $\sigma \notin M_{x}$ then if $\sigma / D=r^{2} / D$ then $\sigma \cdot r^{-2} \in D \subseteq M_{x}$ so
$$
\sigma=\left(\sigma \cdot r^{-2}\right) \cdot r^{2} \in D \cdot M_{x} \subseteq M_{x} \cdot M_{x} \subseteq M_{x}
$$
a contradiction; similarly $\sigma / D \neq 1 / D$. So we have $\sigma / D \in\left\{s / D, s r / D, s r^{2} / D, s r^{3} / D\right\}$. Now, as $M_{x . y} / D \cong \mathbb{Z}_{4}$ we have $M_{x . y}=1 D \cup r^{2} D \cup r D \cup r^{3} D$ (see the proof of Theorem 7.3.8(ii)) and
$$
\left\{M_{x}, M_{y}\right\}=\left\{1 D \cup r^{2} D \cup s D \cup s r^{2} D, 1 D \cup r^{2} D \cup s r D \cup s r^{3} D\right\}
$$

If $M_{x}=1 D \cup r^{2} D \cup s D \cup s r^{2} D$ then as $\sigma \notin M_{x}$ we have $\sigma / D \notin\left\{s / D, s r^{2} / D\right\}$ so

$$
\sigma / D \in\left\{s r / D, s r^{3} / D\right\} \subseteq M_{y} / D
$$

then $\sigma \in M_{y}$ that is $y \in\{\Phi(\mathcal{G}), \sigma \Phi(\mathcal{G})\}^{\perp}$; similarly if $M_{x}=1 D \cup r^{2} D \cup s r D \cup s r^{3} D$ then $y \in\{\Phi(\mathcal{G}), \sigma \Phi(\mathcal{G})\}^{\perp}$.
ii- $\quad$ Let $\bar{T}=\{\Phi(\mathcal{G}), \sigma \Phi(\mathcal{G})\}$.

Claim: To have that $\bar{T}^{\perp} \subseteq G$ is not a saturated subgroup it is enough to prove the following: $\exists x, c \in G \backslash \bar{T}^{\perp}$ such that $<x, c>\equiv<1, x c>$.

Proof of Claim: If $-1 \in \bar{T}^{\perp}$ then $\bar{T}^{\perp} \subsetneq G$ so $G=D_{G}<1,-1>\nsubseteq \bar{T}^{\perp}$ so $\bar{T}^{\perp}$ is not a saturated subgroup. If $-1 \notin \bar{T}^{\perp}$ then take $x, c \in G \backslash \bar{T}^{\perp}$ such that $<x, c>\equiv<1, x c>$ so we have $<c,-x c>\equiv<1,-x>$, that is $c \in D_{G}<1,-x>$ and $-x \in \bar{T}^{\perp}$ : if $-x \notin \bar{T}^{\perp}$ then as $\bar{T}^{\perp} \subseteq G$ has index $2-1 . \bar{T}^{\perp}=-x . \bar{T}^{\perp}$ so $x=-1 .-x \in \bar{T}^{\perp}$; that is we established that there are $a(=-x) \in \bar{T}^{\perp}$ and $c \in D_{G}<1, a>$ with $c \notin \bar{T}^{\perp}$ : this means that $\bar{T}^{\perp}$ is not saturated. Now we will prove that $\sigma^{2} \neq i d$ entails $\exists x, c \in G \backslash \bar{T}^{\perp}$ such that $<x, c>\equiv<1, x c>$ : Take $B$ any well ordered base of $G$ and consider the composition

$$
\mathcal{W}(B) \rightarrow \mathcal{W}(B) / \mathcal{V}(B) \stackrel{\cong}{\rightrightarrows} \mathcal{G}
$$

and, by Lemma 7.1.10, choose any lifting $\widetilde{\sigma} \in \mathcal{W}(B)$ of $\sigma \in \mathcal{G}$. Since $\sigma^{2} \neq i d \in \mathcal{G}$ we get
$\tilde{\sigma}^{2} \in \Phi(B) \backslash \mathcal{V}(B)$. Since

$$
\begin{aligned}
\mathcal{V}(B) & =Q_{B}^{\perp}=\left[\left\{q_{a b}^{B} \in P_{2}(B): l(a) l(b)=0 \in k_{2}(G)\right\}\right]^{\perp} \\
& =\bigcap\left\{\left(q_{a b}^{B}\right)^{\perp} \subseteq \Phi(B): a, b \neq 1, l(a) l(b)=0\right\}
\end{aligned}
$$

we get $a, b \neq 1$ with $l(a) l(b)=0 \in k_{2}(G)$ and $\widetilde{\sigma}^{2} \in \Phi(B) \backslash\left(q_{a b}^{B}\right)^{\perp}$. There are two cases to consider: $a=b$ and $a \neq b$. In the first case $\{a\}$ is a singleton l.i. set and in the second $\{a, b\}$ is a two element l.i. set: consider any well ordered basis $B^{\prime}=\left\{a_{i}^{\prime}: i \in I\right\}$ such that $a=a_{i}$ for some $i \in I$ in the first case and, $a=a_{i}, b=a_{j}$ for some $i<j \in I$ in the second case. Now consider the isomorphism of change of basis $\mu_{B^{\prime} B}: \mathcal{W}(B) \xlongequal{\cong} \mathcal{W}\left(B^{\prime}\right)$ (see Lemma 7.2.3) and take $\sigma^{\prime}=\mu_{B^{\prime} B}(\widetilde{\sigma}) \in \mathcal{W}\left(B^{\prime}\right)$. Then $\sigma^{\prime} \in \mathcal{W}\left(B^{\prime}\right)$ is a lifting of $\sigma \in \mathcal{G}$ with respect to the epimorphism $\mathcal{W}\left(B^{\prime}\right) \rightarrow \mathcal{W}\left(B^{\prime}\right) / \mathcal{V}\left(B^{\prime}\right) \xlongequal{\cong} \mathcal{G}$ (Lemma 7.1.10, again) and $\sigma^{\prime 2} \in \Phi\left(B^{\prime}\right) \backslash \mu_{B^{\prime} B}\left[\left(q_{a, b}^{B}\right)^{\perp}\right]=\Phi\left(B^{\prime}\right) \backslash\left(q_{a, b}^{B^{\prime}}\right)^{\perp}$ (see Remark 7.2.4). As $l(a) l(b)=0$, by the proof of the Theorem 7.3.8, we have $\left(q_{a_{i}, a_{i}}^{B}\right)^{\perp}=S_{i}^{\prime} \cap \Phi\left(B^{\prime}\right)$ in first case and $\left(q_{a_{i}, a_{j}}^{B}\right)^{\perp}=D_{i j}^{\prime} \cap \Phi\left(B^{\prime}\right)$ in the second case, then $\sigma^{\prime 2} \in \Phi\left(B^{\prime}\right) \backslash S_{i}^{\prime}$ (resp. $\sigma^{\prime 2} \in \Phi\left(B^{\prime}\right) \backslash D_{i j}^{\prime}$ ). A straightforward calculation with the group operation in $\mathcal{W}\left(B^{\prime}\right)$ gives $\sigma^{\prime} \notin M_{i}^{\prime} \subseteq \mathcal{W}\left(B^{\prime}\right)$ in the first case and $\sigma^{\prime} \notin M_{i}^{\prime} \cup M_{j}^{\prime} \subseteq \mathcal{W}\left(B^{\prime}\right)$ in the second case, then applying $\mathcal{W}\left(B^{\prime}\right) \rightarrow \mathcal{G}$ we have $\sigma \notin M_{a_{i}} \subseteq \mathcal{G}$ (resp. $\sigma \notin M_{a_{i}} \cup M_{a_{j}} \subseteq \mathcal{G}$ ). Now recall that for each $y \in G \backslash\{1\}$ and each $\theta \in \mathcal{G} \backslash \Phi(\mathcal{G})$, $\theta \notin M_{y}$ iff $<\{\Phi(\mathcal{G}), \theta \cdot \Phi(\mathcal{G})\}, y>=1$ iff $y \notin\{\Phi(\mathcal{G}), \theta \cdot \Phi(\mathcal{G})\}^{\perp}$. Then, since $G$ is a standard prespecial group we have, in both cases, $1 \neq a, b, l(a) l(b)=0 \in k_{2}(G), a, b \notin\{\Phi(\mathcal{G}), \sigma . \Phi(\mathcal{G})\}^{\perp}$ and, in particular, since $G$ is a $k$-stable pre-special group, $1 \in D_{G}(\langle a, b\rangle)$ or, equivalently, $\langle a, b\rangle \equiv\langle 1, a b\rangle$.
iii- Recall that for special groups the maximal saturated subgroups are precisely the index 2 saturated subgroups so the result follows from items (i) and (ii).
iv- Let $\Delta \subseteq G$ be a saturated subgroup: as $G$ is a special group $\Delta=\bigcap\left\{\Sigma \subseteq G: \Sigma \in X_{\Delta}\right\}$ where

$$
X_{\Delta}=\{\Sigma \subseteq G: \Sigma \text { is a maximal saturated subgroup and } \Delta \subseteq \Sigma\}
$$

then, by Proposition 7.3.4,

$$
\Delta^{\perp}=\bigvee\left\{\Sigma^{\perp} \subseteq \mathcal{G} / \Phi(\mathcal{G}): \Sigma \in X_{\Delta}\right\}
$$

by item (iii) $\Sigma^{\perp}=\{\Phi(\mathcal{G}), \sigma \Phi(\mathcal{G})\}$ for some $\sigma \in \mathcal{G} \backslash \Phi(\mathcal{G}), \sigma^{2}=i d$; take $T_{\Sigma}=\Phi(\mathcal{G}) \cup \sigma \Phi(\mathcal{G})$ then $T_{\Sigma} \subseteq \mathcal{G}$ is a closed (normal) subgroup such that $\Phi(\mathcal{G}) \subseteq T_{\Sigma}$ and all elements of $T_{\Sigma} \backslash\{1\}$ are involutions so

$$
\bigvee\left\{T_{\Sigma}: \Sigma \in X_{\Delta}\right\}=\operatorname{closure}\left(\left[\left\{T_{\Sigma}: \Sigma \in X_{\Delta}\right\}\right]\right)
$$

is a closed subgroup of $\mathcal{G}$ that contains $\Phi(\mathcal{G})$ and is (topologically) generated by involutions. Now note that $T_{\Sigma} / \Phi(\mathcal{G})=\{\Phi(\mathcal{G}), \sigma \Phi(\mathcal{G})\}=\Sigma^{\perp}$ and then

$$
\begin{aligned}
\Delta^{\perp} & =\bigvee\left\{\Sigma^{\perp} \subseteq \mathcal{G} / \Phi(\mathcal{G}): \Sigma \in X_{\Delta}\right\}=\bigvee\left\{T_{\Sigma} / \Phi(\mathcal{G}) \subseteq \mathcal{G} / \Phi(\mathcal{G}): \Sigma \in X_{\Delta}\right\} \\
& =\left(\bigvee\left\{T_{\Sigma} \subseteq \mathcal{G}: \Sigma \in X_{\Delta}\right\}\right) / \Phi(\mathcal{G})
\end{aligned}
$$

as $q: \mathcal{G} \rightarrow \mathcal{G} / \Phi(\mathcal{G})$ gives an isomorphism of complete lattices between the set of closed subgroups of $\mathcal{G}$ which contains $\Phi(\mathcal{G})$ and the set of closed subgroups of $\mathcal{G} / \Phi(\mathcal{G})$.
Now take $T \subseteq \mathcal{G}$ a closed subgroup of $\mathcal{G}$ such that $\Phi(\mathcal{G}) \subseteq T$ and $T$ is topologically generated
by involutions. Write $I_{T}=\left\{\sigma \in T: \sigma \in \mathcal{G} \backslash \Phi(\mathcal{G}), \sigma^{2}=i d\right\}$ then, for each $\sigma \in I_{T}, \sigma \Phi(\mathcal{G}) \subseteq T$ , $T_{\sigma}=\Phi(\mathcal{G}) \cup \sigma \Phi(\mathcal{G})$ is a closed (normal) subgroup of $\mathcal{G}$ and

$$
T=\operatorname{closure}\left(\left\lfloor\left\{T_{\sigma}: \sigma \in I_{T}\right\}\right]\right)=\bigvee\left\{T_{\sigma}: \sigma \in I_{T}\right\}
$$

also $T_{\sigma} / \Phi(\mathcal{G})=\{\Phi(\mathcal{G}), \sigma \Phi(\mathcal{G})\}$ and, by item (iv), $\left(T_{\sigma} / \Phi(\mathcal{G})\right)^{\perp} \subseteq G$ is a maximal saturated subgroup of $G$. Then we have

$$
\begin{aligned}
(T / \Phi(\mathcal{G}))^{\perp} & =\left(\left(\bigvee\left\{T_{\sigma}: \sigma \in I_{T}\right) / \Phi(\mathcal{G})\right)^{\perp}=\left(\bigvee\left\{T_{\sigma} / \Phi(\mathcal{G}): \sigma \in I_{T}\right\}\right)^{\perp}\right. \\
& =\bigcap\left\{\left(T_{\sigma} / \Phi(\mathcal{G})\right)^{\perp}: \sigma \in I_{T}\right\}
\end{aligned}
$$

which is a saturated subgroup of $G$.

Theorem 7.3.13. Let $G$ be a standard special group. Are equivalent
$i$ - G is "Pythagorean" or "almost reduced, 2
ii - $\operatorname{Gal}(G)$ is generated by involutions.
Proof. Note that the unique non-formally real Pythagorean special group (equivalently, $-1 \neq 1$ ) is $G=\{1\}$ and thus $\operatorname{Gal}(G)=\{1\}$.
$(i) \Rightarrow(i i)$ : The hypothesis means that $\{1\} \subseteq G$ is a saturated subset of $G$, then by item (iv) of the previous Proposition, $\operatorname{Gal}(G)$ is generated by involutions.
$(i i) \Rightarrow(i)$ : It follows the hypothesis that there is no continuous epimorphism $\mathcal{G} / \Phi(\mathcal{G}) \rightarrow \mathbb{Z}_{4}$. Since $G$ is standard SG, for all $a \in G \backslash\{1\}, l(a) l(a) \neq 0=l(1) l(1)$ and, since $G$ is in particular $k$-stable, then for all $a \in G \backslash\{1\}$, is not the case $\langle a, a\rangle \equiv<1,1\rangle$, that is: $G$ is Pythagorean.

Remark 7.3.14. Another Galois theoretic characterization of the Pythagoreaness of $G$ is

$$
\Phi(\operatorname{Gal}(G))=[\operatorname{Gal}(G), \operatorname{Gal}(G)] .
$$

Theorem 7.3.15. Let $G$ be a standard special group. Consider the following
$i$ - $G$ is not formally real.
ii- Every involution is is $\Phi(\operatorname{Gal}(G))$.
iii- Every involution in $\operatorname{Gal}(G)$ is central.
Then $(i) \Rightarrow(i i) \Rightarrow$ (iii) and if $\operatorname{card}(G a l(G))>2$, then all are equivalent.
Proof. By Theorem 7.3.12(iii) $G$ is formally real iff there is an involution $\sigma \in \operatorname{Gal}(G) \backslash \Phi(\operatorname{Gal}(G))$, so we get $(i) \Rightarrow(i i)$. As $\operatorname{Gal}(G)$ is a $\mathcal{C}$-group we have $\left[\sigma^{2}, \tau\right]=1$ and, since $\Phi(\operatorname{Gal}(G))=\operatorname{Gal}(G)^{2}$ (pro-2-group), then $\Phi(\operatorname{Gal}(G)) \subseteq \operatorname{center}(\operatorname{Gal}(G))$ so $(i i) \Rightarrow(i i i)$. Now suppose $\operatorname{card}(\operatorname{Gal}(G))>2$ : to prove $(i i i) \Rightarrow(i)$ let us assume $G$ formally real and note that any involution $\notin \Phi(G a l(G))$ is not in the center of $\operatorname{Gal}(G)$.

[^17]
### 7.4 The functorial behavior of $G a l$ and SG-cohomology

We have developed the theme "basis change induced isomorphisms" in the general context of pre-Special Groups, as the fundamental step to get a single Galois group of a pre-special group. In this final section, we analyze some functorial behavior of the Gal construction of $S G$-theory and provide the first steps to a (profinite) "Galoisian" cohomology for the $S G$-theory, in an attempt to complete the "Milnor scenario" of Igr's (52]) in abstract theories of quadratic forms.

### 7.4.1 From PSG to Galois groups

### 7.4.1. Construction:

Let $f: G \rightarrow G^{\prime}$ a pSG-homomorphism of pre-special groups.
Let $B_{1}^{\prime}=\left\{a_{k}^{\prime}: k \in I_{1}^{\prime}\right\}$ be an well ordered basis of $f[G]$ and extends it to $B^{\prime}=\left\{a_{k}^{\prime}: k \in I^{\prime}\right\}$, an well ordered basis of $G^{\prime}$. Now select $a_{k} \in f^{-1}\left[\left\{a_{k}^{\prime}\right\}\right], k \in I_{1}^{\prime}$. Then the set $B_{1}=\left\{a_{i}: i \in I_{1}^{\prime}\right\} \subseteq G$ is linearly independent, now complete this to basis of $G, B=\left\{a_{i}: i \in I\right\}$ : we just need to glue a well ordered basis of $\operatorname{ker}(f)$.

We have some induced functions:
(0) $f_{B, B^{\prime}}^{0}: B^{\prime} \rightarrow \mathcal{W}(B)$ is such that $f_{B, B^{\prime}}^{0}\left(a_{k}^{\prime}\right)=a_{k}$, if $a_{k}^{\prime} \in B_{1}^{\prime}$ and $f_{B, B^{\prime}}^{0}\left(a_{k}^{\prime}\right)=1$, if $a_{k}^{\prime} \in B^{\prime} \backslash B_{1}^{\prime}$.
(1) $f_{B, B^{\prime}}^{(1)}: B \rightarrow P_{2}\left(B^{\prime}\right)$ is such that $f_{B, B^{\prime}}^{1}\left(a_{k}\right)=a_{k}^{\prime}$, if $a_{k} \in B_{1}$ and $f_{B, B^{\prime}}^{1}\left(a_{k}\right)=0$, if $a_{k} \in B \backslash B_{1}$.
(2) $f_{B, B^{\prime}}^{(2)}:\left\{\left(a_{i}, a_{j}\right) \in B \times B: i, j \in I, i \leq j\right\} \rightarrow P_{2}\left(B^{\prime}\right)$ is such that $f_{B, B^{\prime}}^{2}\left(a_{i}, a_{j}\right)=a_{i}^{\prime} \cdot a_{j}^{\prime}$, if $\left(a_{i}, a_{j}\right) \in B_{1} \times B_{1}$ and $f_{B, B^{\prime}}^{2}\left(a_{i}, a_{j}\right)=0$, if $\left(a_{i}, a_{j}\right) \in B \times B \backslash B_{1} \times B_{1}$.

Keeping the notation above, we have
Proposition 7.4.2. The function $f_{B, B^{\prime}}^{0}: B^{\prime} \rightarrow \mathcal{W}(B)$ induces a continuous homomorphism $\bar{f}_{B, B^{\prime}}$ : $\mathcal{W}\left(B^{\prime}\right) / \mathcal{V}\left(B^{\prime}\right) \rightarrow \mathcal{W}(B) / \mathcal{V}(B)$.

Proof. It follows from Proposition 7.1.6 and the definition of $f_{B, B^{\prime}}^{0}: B^{\prime} \rightarrow \mathcal{W}(B)$, that its image converges to $1 \in \mathcal{W}(B)$. Thus, by the universal property $\mathcal{W}\left(B^{\prime}\right)$ (Theorem 7.1.11), $f_{B, B^{\prime}}^{0}$ extends uniquely to a continuous homomorphism of pro-2-groups $\hat{f}_{B, B^{\prime}}^{0}: \mathcal{W}\left(B^{\prime}\right) \rightarrow \mathcal{W}(B)$.

Now, $f: G \rightarrow G^{\prime}$ also induces a $\mathbb{Z}_{2}$-module homomorphism $\hat{f}_{B, B^{\prime}}^{(2)}: P_{2}(B) \rightarrow P_{2}\left(B^{\prime}\right)$ : this is just the unique $\mathbb{Z}_{2}$-linear extension of the induced map $f_{B, B^{\prime}}^{(2)}:\left\{\left(a_{i}, a_{j}\right) \in B \times B: i, j \in I, i \leq\right.$ $j\} \rightarrow P_{2}\left(B^{\prime}\right)$. Moreover, since $k_{*}(f): k_{*}(G) \rightarrow k_{*}\left(G^{\prime}\right)$ is an Igr-morphism, then for each $a, b \in G$ such that $l(a) l(b)=0 \in k_{2}(G)$, we have $l(f a) \cdot l(f b)=0 \in k_{2}\left(G^{\prime}\right)$. From this we obtain that $\hat{f}_{B, B^{\prime}}^{(2)}\left(q_{a, b}^{B}\right)=q_{f a, f b}^{B^{\prime}}$ (see Proposition 7.2 .2 and, therefore, $f_{B, B^{\prime}}^{(2)}[Q(B)] \subseteq Q\left(B^{\prime}\right)$.

For each $a, b \in G$ and $\sigma^{\prime} \in \mathcal{W}(B)$, we have

$$
<\hat{f}_{B, B^{\prime}}^{(0)}\left(\sigma^{\prime}\right), q_{a, b}^{B}>=<\sigma^{\prime}, \hat{f}_{B, B^{\prime}}^{(2)}\left(q_{a, b}^{B}>=<\sigma^{\prime}, q_{f a, f b}^{B^{\prime}}>\right.
$$

Thus $\hat{f}_{B, B^{\prime}}^{(0)}\left[\mathcal{V}\left(B^{\prime}\right)\right] \subseteq \mathcal{V}(B)$ and then, $\hat{f}_{B, B^{\prime}}^{0}: \mathcal{W}\left(B^{\prime}\right) \rightarrow \mathcal{W}(B)$ induces a unique continuous homomorphism of pro-2-groups $\bar{f}_{B, B^{\prime}}: \mathcal{W}\left(B^{\prime}\right) / \mathcal{V}\left(B^{\prime}\right) \rightarrow \mathcal{W}(B) / \mathcal{V}(B)$.

Proposition 7.4.3. Let $f: G \rightarrow G^{\prime}$ be an injective $p S G$-morphism.
$i$ - Then $\hat{f}_{B, B^{\prime}}^{0}: \mathcal{W}\left(B^{\prime}\right) \rightarrow \mathcal{W}(B)$ and $\bar{f}_{B, B^{\prime}}: \mathcal{W}\left(B^{\prime}\right) / \mathcal{V}\left(B^{\prime}\right) \rightarrow \mathcal{W}(B) / \mathcal{V}(B)$ are surjective morphisms of pro-2-groups (thus they can be identified with projections).
ii - Let $f^{\prime}: G^{\prime \prime} \rightarrow G$ is an injective $p S G$-morphism. If $B^{\prime}$ is an well ordered basis of $G^{\prime}$ obtained by successive extensions of an well ordered basis $B_{2}^{\prime}$ of $f \circ f^{\prime}\left[G^{\prime \prime}\right]$ to $B_{1}^{\prime}$, an well ordered basis of $f[G] \supseteq f \circ f^{\prime}\left[G^{\prime \prime}\right]$, then applying the construction above described, we obtain
$\hat{f}_{B^{\prime \prime}, B^{\prime}}^{0}=\hat{f}_{B^{\prime \prime}, B}^{0} \circ \hat{f}_{B, B^{\prime}}^{0}: \mathcal{W}\left(B^{\prime}\right) \rightarrow \mathcal{W}\left(B^{\prime \prime}\right)$ and $\overline{f_{B^{\prime \prime}, B^{\prime}}^{0}}=\bar{f}_{B^{\prime \prime}, B}^{0} \circ \overline{f_{B, B^{\prime}}^{0}}: \mathcal{W}\left(B^{\prime}\right) / \mathcal{V}\left(B^{\prime}\right) \rightarrow$ $\mathcal{W}\left(B^{\prime \prime}\right) / \mathcal{V}\left(B^{\prime \prime}\right)$

## Proof.

i - By the injectivity hypothesis, we have $f_{B, B^{\prime}}^{0}\left[B^{\prime}\right]=B \cup\{1\} \subseteq \mathcal{W}(B)$, thus $\hat{f}_{B, B^{\prime}}^{0}: \mathcal{W}\left(B^{\prime}\right) \rightarrow$ $\mathcal{W}(B)$ is a continuous function with dense image from a compact space into a Hausdorff space. Therefore $\hat{f}_{B, B^{\prime}}^{0}$ and $\bar{f}_{B, B^{\prime}}^{0}$ are surjective continuous homomorphisms.
ii - It follows from a straightforward calculation that $f_{B^{\prime \prime}, B^{\prime}}^{0}=\hat{f}_{B^{\prime \prime}, B}^{0} \circ f_{B, B^{\prime}}^{0}$. Therefore, the uniqueness of extensions and the homomorphism theorem guarantees that $\hat{f}_{B^{\prime \prime}, B^{\prime}}^{0}=\hat{f}_{B^{\prime \prime}, B}^{0} \circ$ $\hat{f}_{B, B^{\prime}}^{0}$ and $\overline{f_{B^{\prime \prime}}^{0}, B^{\prime}}=\bar{f}_{B^{\prime \prime}, B}^{0} \circ \bar{f}_{B, B^{\prime}}^{0}$.

### 7.4.2 From Galois Groups to PSG

Let $G$ be a pre-special group an denote $\mathcal{G}=\operatorname{Gal}(G)$. We have seen in Proposition 7.3.2 that there is a canonical isomorphism $\phi_{G}: \mathcal{G} / \Phi(\mathcal{G}) \stackrel{\cong}{\rightrightarrows} \operatorname{Hom}\left(G, \mathbb{Z}_{2}\right)$ so we get a "perfect pairing" $\hat{\phi}_{G}: \mathcal{G} / \Phi(\mathcal{G}) \times G \rightarrow \mathbb{Z}_{2}$ and there is also a canonical bijection $G \cong\{T \subseteq \mathcal{G}: T$ is a closed normal subgroup of index $\leq 2\}$.

We will explain now the term "canonical" employed, starting with the following

## Lemma 7.4.4.

$i$ - Let $G, G^{\prime}$ be pre-special groups Then each continuous homomorphism $\theta: \operatorname{Gal}\left(G^{\prime}\right) \rightarrow \operatorname{Gal}(G)$ induces a $\mathbb{Z}_{2}$-module homomorphism $\check{\theta}: G \rightarrow G^{\prime}$.
ii - The association above, $\theta \mapsto \check{\theta}$, determines a contravariant functor from the category from all pairs $(G, \operatorname{Gal}(G)), G$ a pre-special group, and continuous homomorphisms, into the category of $\mathbb{Z}_{2}$-modules.

## Proof.

i - We have a $\mathbb{Z}_{2}$-homomorphism $\theta^{*}: \operatorname{Homcont}\left(\operatorname{Gal}(G), \mathbb{Z}_{2}\right) \rightarrow \operatorname{Homcont}\left(\operatorname{Gal}\left(G^{\prime}\right), \mathbb{Z}_{2}\right), \mu \mapsto \mu \circ \theta$. By Proposition 7.3 .2 (iii), we have $\mathbb{Z}_{2}$-isomorphisms $\psi_{G}, \psi_{G^{\prime}}$. Combining the informations we define the $\mathbb{Z}_{2}$-homomorphism $\check{\theta}:=\psi_{G^{\prime}}^{-1} \circ \theta^{*} \circ \psi_{G}: G \rightarrow G^{\prime}$.
ii - Note that $i d_{\mathcal{G}}^{*}=i d$, thus $\check{i d} d_{\mathcal{G}}=i d_{G}$. Let $\theta^{\prime}: \operatorname{Gal}\left(G^{\prime \prime}\right) \rightarrow \operatorname{Gal}\left(G^{\prime}\right)$ be a continuous homomorphism. Then we have $\left(\theta \circ \theta^{\prime}\right)^{*}=\theta^{\prime *} \circ \theta^{*}$, thus $\left(\theta \circ \theta^{\prime}\right)^{\prime}=\check{\theta}^{\prime} \circ \check{\theta}$.

Remark 7.4.5. Note that for any surjective continuous homomorphism $\theta: \mathcal{G}^{\prime} \rightarrow \mathcal{G}$, we have that $\check{\theta}: G \rightarrow G^{\prime}$ is an injective $\mathbb{Z}_{2}$-homomorphism.

This suggest that the (sub)category of Galois groups and continuous epimorphisms is the "right" domain category of Galois groups. We have the following:

Proposition 7.4.6. The functor described in the Lemma above restricts to a functor from the subcategory of Galois groups of standard pre-special groups (Definition above 7.3.9) and continuous epimorphisms to the category of standard pre-special groups and injective $q S G$-morphisms (i.e., the group homomorphisms that preserves $\equiv$, but that eventually does not preserves -1).

Proof. Assume that $G$ is a $k$-stable pre-special group and that $G^{\prime}$ is a standard pre-special group, we will prove that $\check{\theta}$ is a injective $q S G$-homomorphism from $G$ to $G^{\prime}$.

Since $\check{\theta}: G \rightarrow G^{\prime}$ is a group homomorphism, it is enough to show that, for each $a, b \in G \backslash\{1\}$, $1 \in D_{G}<a, b>\Rightarrow 1^{\prime} \in D_{G}<\check{\theta}(a), \check{\theta}(b)>$ and, since $G, G^{\prime}$ are $k$-stable pre-special group, this is equivalent to show $l(a) l(b)=0 \in k_{2}(G) \Rightarrow l(\check{\theta}(a)) l(\check{\theta}(b))=0 \in k_{2}\left(G^{\prime}\right)$. Now we have 2 cases to consider:
$\underline{a=b}$ : Since $G$ is a $k$-stable pre-special group then, by Theorem 7.3.8(i), there exists $S \subseteq \mathcal{G}$ a closed normal subgroup with $\mathcal{G} / S \cong \mathbb{Z}_{4}$ and $S \subseteq M_{a}$. As $\theta: \mathcal{G}^{\prime} \rightarrow \mathcal{G}$ is an surjective continuous homomorphism we have that the quotient map $\theta_{S}: \mathcal{G}^{\prime} / \theta^{-1}[S] \rightarrow \mathcal{G} / S$ is an isomorphism so $\theta^{-1}[S] \subseteq$ $\mathcal{G}^{\prime}$ is a closed normal subgroup such that $\mathcal{G}^{\prime} / \theta^{-1}[S] \cong \mathcal{G} / S \cong \mathbb{Z}_{4}$ and $\theta^{-1}[S] \subseteq \theta^{-1}\left[M_{a}\right]=M_{\theta^{\star}(a)}^{\prime}$, where this last equality holds by the bijections in items (ii) and (iii) in Proposition 7.3.2 and the definition of $\check{\theta}$. Now, since $G^{\prime}$ is a standard pre-special group, we have $l(\check{\theta}(a)) l(\check{\theta}(a))=0 \in k_{2}\left(G^{\prime}\right)$
$a \neq b:$ As $\check{\theta}$ is an injective $\mathbb{Z}_{2}$-homomorphism, $\check{\theta}(a) \neq \check{\theta}(b)$. Since $G$ is a $k$-stable pre-special group then, by Theorem 7.3 .8 (ii), there exists $D \subseteq \mathcal{G}$ a closed normal subgroup with $\mathcal{G} / D \cong \mathbb{D}_{4}$, $D \subseteq M_{a} \cap M_{b}$ and $M_{a b} / D \cong \mathbb{Z}_{4}$. As $\theta: \mathcal{G}^{\prime} \rightarrow \mathcal{G}$ is an surjective continuous homomorphism we have that the quotient map $\theta_{D}: \mathcal{G}^{\prime} / \theta^{-1}[D] \rightarrow \mathcal{G} / D$ is an isomorphism so $\theta^{-1}[D] \subseteq \mathcal{G}^{\prime}$ is a closed normal subgroup such that $\mathcal{G}^{\prime} / \theta^{-1}[D] \cong \mathcal{G} / D \cong \mathbb{D}_{4}$ and $\theta^{-1}[D] \subseteq \theta^{-1}\left[M_{a} \cap M_{b}\right]=M_{\hat{\theta}(a)}^{\prime} \cap M_{\tilde{\theta}(b)}^{\prime}$. Now we check that $M_{\tilde{\theta}(a) \tilde{\theta}(b)}^{\prime} / \theta^{-1}[D] \cong \mathbb{Z}_{4}$ : as $\theta$ is an epimorphism $\theta\left[\right.$ theta $\left.^{-1}\left[M_{a b}\right]\right]=M_{a b}$ so, as $M_{\ddot{\theta}(a) \check{\theta}(b)}^{\prime}=M_{\stackrel{\theta}{(a b)}}^{\prime}=\theta^{-1}\left[M_{a b}\right]$ (because $\check{\theta}(a) \neq \check{\theta}(b)$ and there are exactly three maximal above $\left.M_{\ddot{\theta}(a)}^{\prime} \cap M_{\grave{\theta}(b)}^{\prime}\right)$, we have an epimorphism $\theta_{\mid}: M_{\grave{\theta}(a b)}^{\prime} \rightarrow M_{a b}$ thus $\operatorname{ker}\left(\theta_{\mid}\right)=\theta^{-1}[D] \subseteq M_{\ddot{\theta}(a b)}^{\prime}$ and the quotient map $\theta_{\mid D}: M_{\hat{\theta}(a b)}^{\prime} / \theta^{-1}[D] \rightarrow M_{a b} / D$ is an isomorphism so $M_{\hat{\theta}(a b)}^{\prime} / \theta^{-1}[D] \cong M_{a b} / D \cong \mathbb{Z}_{4}$. Now, since $G^{\prime}$ is a standard pre-special group, we have $l(\check{\theta}(a)) l(\check{\theta}(b))=0 \in k_{2}\left(G^{\prime}\right)$.

### 7.4.3 Towards a galoisian cohomology for $S G$-theory

Let $G$ be a standard pre-special group. Since $\mathcal{G}=\operatorname{Gal}(G)$ is a profinite group, the Galois Cohomology is available for this subclass of pre-special groups. In particular, there is the graded cohomology ring $H^{*}(G):=H^{*}(\mathcal{G},\{ \pm 1\})$ where $\mathcal{G}$ act trivially on $\mathbb{F}_{2}=\mathbb{Z} / 2 \mathbb{Z}$. Therefore, at least some parts of Milnor's scenario for containing 3 graded rings related to quadratic forms theory of fields (with char $\neq 2$ ) is available for (standard) special groups: $W_{*}(G), k_{*}(G)$, and $H^{*}(G)$.

The result above just provides the initial step to establish cohomological methods in SG-theory.
Theorem 7.4.7. As in the field case, consider $\mathbb{Z}_{2} \cong\{ \pm 1\}$ as a discrete $\operatorname{Gal}(G)$-module endowed with the trivial action, i.e., $\sigma . a=a$, for all $\sigma \in \operatorname{Gal}(G)$ and $a \in \mathbb{Z}_{2}$. Then $H_{*}(G):=$ $H_{*}\left(\operatorname{Gal}(G), \mathbb{Z}_{2}\right)$, is an Igr, endowed with the cup product. Moreover, there is a canonical isomorphism of pointed 2 -groups $(G,-1) \cong\left(H^{1}(G),(-1)\right)$.

Proof. We write $\mathcal{G}:=\operatorname{Gal}(G)$. Just recall that:
$H^{0}\left(\mathcal{G}, \mathbb{Z}_{2}\right)=\left(\mathbb{Z}_{2}\right)^{\mathcal{G}}=\operatorname{Fix}\left(\mathbb{Z}_{2}\right)=\mathbb{Z}_{2}$, since $\mathcal{G}$ is acting trivially on $\mathbb{Z}_{2}$.
For $H^{1}\left(\mathcal{G}, \mathbb{Z}_{2}\right):=\operatorname{Crossed} \operatorname{Hom}\left(\mathcal{G}, \mathbb{Z}_{2}\right) /$ principalCrossed $\operatorname{Hom}\left(\mathcal{G}, \mathbb{Z}_{2}\right)$, since $\mathcal{G}$ is acting trivially
on $\mathbb{Z}_{2}$, we get

$$
\begin{aligned}
\text { principalCrossedHom }\left(\mathcal{G}, \mathbb{Z}_{2}\right) & :=\operatorname{Im}\left(\bar{\partial}_{1}\right) \\
& :=\left\{x: \mathcal{G} \rightarrow \mathbb{Z}_{2}: x=\bar{\partial}_{1} a \text { for some } a \in \mathbb{F}_{2}\right\} \\
& =\left\{x: \mathcal{G} \rightarrow \mathbb{Z}_{2}: \text { there exist } \mathbb{F}_{2} \in \mathbb{F}_{2} \text { such that } x(\sigma)=\sigma a-a \text { for all } \sigma \in \mathcal{G}\right\} \\
& =\left\{x: \mathcal{G} \rightarrow \mathbb{Z}_{2}: \text { there exist } a \in \mathbb{F}_{2} \text { such that } x(\sigma)=0 \text { for all } \sigma \in \mathcal{G}\right\} \\
& =\{0\} ;
\end{aligned}
$$

and

$$
\begin{aligned}
\operatorname{CrossedHom}\left(\mathcal{G}, \mathbb{Z}_{2}\right) & :=\operatorname{Ker}\left(\bar{\partial}_{2}\right) \\
& :=\left\{x: \mathcal{G} \rightarrow \mathbb{Z}_{2}: x \text { is continuous and } x(\sigma \tau)=\sigma x(\tau)+x(\sigma) \text { for all } \sigma, \tau \in \mathcal{G}\right\} \\
& =\left\{x: \mathcal{G} \rightarrow \mathbb{Z}_{2}: x \text { is continuous and } x(\sigma \tau)=x(\tau)+x(\sigma) \text { for all } \sigma, \tau \in \mathcal{G}\right\} \\
& =\operatorname{Homcont}\left(\mathcal{G}, \mathbb{Z}_{2}\right) .
\end{aligned}
$$

Therefore $H^{1}\left(\mathcal{G}, \mathbb{Z}_{2}\right)=H \operatorname{Homcont}\left(\mathcal{G}, \mathbb{Z}_{2}\right) /\{0\} \cong \operatorname{Homcont}\left(\mathcal{G}, \mathbb{Z}_{2}\right)=\operatorname{Homcont}\left(\operatorname{Gal}(G), \mathbb{Z}_{2}\right)$.
On the other hand, by Proposition 7.3 .2 (iii) $\psi_{G}: G \stackrel{\cong}{\rightrightarrows} \operatorname{Homcont}\left(G a l(G), \mathbb{Z}_{2}\right)$ as $\mathbb{Z}_{2}$-modules and, $-1 \in G$ corresponds to a open subgroup of $\operatorname{Gal}(G)$ with index $\leq 2$, that corresponds to ( -1 ) in $\operatorname{Homcont}\left(\operatorname{Gal}(G), \mathbb{Z}_{2}\right)$

It is natural ask if the $\operatorname{Igr} H^{*}(\operatorname{Gal}(G),\{ \pm 1\})$ is in the subcategory $\operatorname{Igr}_{h}$ : this depends of an analysis and more explicit description of $H^{2}(\operatorname{Gal}(G),\{ \pm 1\})$. In particular, we will need to analyze the relationship between the equations $l(a) l(b)=0 \in k_{2}(G)$ and $(a) \cup(b)=0 \in H^{2}\left(G a l(G), \mathbb{Z}_{2}\right)$, for $a, b \in G$, in a standard pre-special group $G$. Related to this question is the existence of a Milnor like canonical arrow from the mod 2 k-theory graduated ring of $G$ to the graduated ring of cohomology of G: $h(G): k_{*}(G) \rightarrow H^{*}(G)$.

These generalizes some results in [2] to the context of (pre)special groups where they prove that the cohomology ring $H^{*}\left(\operatorname{Gal}\left(F^{(3)} \mid F\right), \mathbb{Z}_{2}\right)$ contains the cohomology ring $H^{*}\left(\operatorname{Gal}\left(F^{s} \mid F\right), \mathbb{Z}_{2}\right)$ as its subring generated by cup products of level 1 elements. Therefore, it could be interesting also analyze the properties of the subIgr generated in level 1 of the $\operatorname{Igr} H^{*}(\operatorname{Gal}(G),\{ \pm 1\})$, that is possibly a member of the subclass $\mathrm{Igr}_{+}$.

## Chapter 8

## Conclusion and Further Works

After all, we return to our initial diagram


In this present work, our main results concerns to the relation between special groups and graded rings, with much contribution of the theory of multirings/hyperfields. In fact, we established the result of Arason-Pfister Hauptsatz for every special group (Theorem 6.3.2), as an application of "multialgebraic methods" here introduced.

In Chapter 7 we started the investigation of the relations between special groups and profinite 2-groups, towards completing the initial diagram above.

After that, we glance at these roads to follow:

1. We intend to analyze further the introduced notions of formally real semigroups, formally real multirings and quadratic multirings.
2. With Example 2.5 .15 as a prototype, specialize the study of quadratic multirings where every element is the product of a non-zero divisor and an idempotent. This could give some hint about the structure of invertible elements in real semigroups, which until today is not known to be a reduced special group in general.
3. In [25] is constructed a von Neumann hull functor from multiring category and that, when restricted to in semi-real rings, it commutes with real semigroup functor. This allows us to obtain some quadratic forms properties of a semi-real ring by looking to its von Neumann regular hull. It would be interesting to determine what kind of property in the von Neumann hull of a quadratic multiring return to the original structure.
4. The definition and analysis of the structure of Witt ring of more general quadratic structures (non only obtained from special groups): this subject have already appeared in Section 4 of Chapter 2, in connection with (37.
5. Extension of the K-theory framework to more general multirings (for example, to VNmultirings) with quadratic flavour.
6. Compare graded K-theory with graded Witt ring for VN-real semigroups as in the field case (Milnor [52]) and special groups (Dickmann-Miraglia [28]).
7. In the hyperfield case, investigate the extension of the concept of Galois group to hyperfields, comparing the Galois cohomology ring and analyse the existence of some canonical arrow from K-theory to this cohomology ring, in an attempt to recover the Milnor's Conjecture available in the classic algebraic quadratic forms context ([57], 20], 21).
8. The next steps in the program of study algebraic extensions of superfields are a development of Galois theory and Galois cohomology theory, envisaging application to other mathematical theories as abstract structures of quadratic forms and real algebraic geometry ([24, [17, [18]): some parts of this program are under development in [16] and [14].
9. In the vein of the previous item, we will pursue, in particular, further developments of the theory of quadratic extensions of hyperfields and superfields, envisaging the description of Galois groups of special hyperfields "from below". We intend apply this description to obtain further information on the graded cohomology ring of a special group and provide a more complete development of cohomological methods in SG-theory, applying this to obtain a possible obstruction for every reduced SG to satisfy Marshall's conjecture.
10. Since the theory of superfields/hyperfields and the abstract theories of quadratic forms of Special Groups [28] and of Real Semigroups [33] are (or can be seen as) first-order theories, we wonder about other possible model-theoretic results in these theories. In connection with this, we plan to develop an order theory of superfields and analyze some candidates for notions of real closed superfields in such a way that we may address the questions: (i) the class of real closed superfields admits quantifiers elimination or is model-complete (according to a convenient choice of language)?; (ii) any reasonably ordered superfield admits an essentially unique real closure?
11. It could be interesting describe and explore an alternative notion of algebraically closed multifield based on an alternative notion of of root of a polynomial, taking in account factorizations, for example, if $p(x) \in(x-b) q(x)$ for some $q(x)$, then $b$ can be seem as a root of $p(x)$ : by Theorem 7 in [6], this in fact coincide with the other notion of root of a polynomial $p(x) \in F[x]$ whenever $F$ is a hyperfield.
12. In [55] was started the development of a identity theory and a universal algebra like theory for multi structures. However, a full model theory of multi structures, in the vein of Chapter 1 of [26], should be an object of interest (as the present work suggests) and it is seems to be unknown.
13. Examples 3.6.16 and 3.7.7 reveals the necessity of some computational implementation in order to ease and accelerate the calculations with algebraic extensions of superfields: in [13], we start a proposal towards this subject.

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[^0]:    ${ }^{1}$ Carried out mainly by I. Minac and co-authors.

[^1]:    ${ }^{1}$ Here the subscript "fol" is to indicate that we are thinking in the first order theory associated to multigroups.

[^2]:    ${ }^{2}$ According Marshall's notation in 47.

[^3]:    ${ }^{3}$ There is no consensus on the definition "submultiring": here we do adopted one of intermediary strength that coincides with the notion of substructure in relational structures; in 47, submultiring means an inclusion of multirings that is strong and full.

[^4]:    ${ }^{1}$ The name "Dickmann-Miraglia" is given in honor to professors Maximo Dickmann and Francisco Miraglia, the creators of the special group theory.

[^5]:    ${ }^{2}$ The name "Dickmann-Petrovich" is given in honor to professors Max Dickmann and Alejandro Petrovich, who are the creators of realsemigroup theory.

[^6]:    ${ }^{1}$ For a quasi-superfield $F$, we are not imposing that $(S \backslash\{0\}, \cdot, 1)$ will be a commutative multigroup, i.e, that if $d \in a \cdot b$ then $b^{-1} \in a \cdot d^{-1}$.

[^7]:    ${ }^{2}$ A superring $B$ is called quasi-superdomain if given $a, b \in B$ with $a b=\{0\}$, then $a=0$ or $b=0$

[^8]:    ${ }^{3}$ Or saying in another words, in order to obtain $a_{1} x_{1}+b_{1} x_{1}=\left(a_{1}+b_{1}\right) x_{1}$ we should Define what would be a "full" variable, which is not a standard procedure in logic.

[^9]:    ${ }^{4}$ Elementary matrices are a standard topic in many Linear Algebra books, but for a quick reference, consult https://en.wikipedia.org/wiki/Elementary_matrix.

[^10]:    ${ }^{5}$ Of course, not necessarily unique.

[^11]:    ${ }^{6}$ As we will see later, simple calculations with superfield are highly demanding...

[^12]:    ${ }^{1}$ The main terminology in the literature is "hyperring". Moreover, M. Marshall makes a distinction between "multiring" and "hyperrings" which is important in the context of quadratic forms. But throughout this entire work, we deal essentially with multifields/hyperfields and then, the main terminology here will be "hyperfield".

[^13]:    ${ }^{2}$ Remember that we are using the simplified notation for elements in $K_{n}^{\text {mult }} F$ (and all other K-theories), which is $\rho\left(a_{1}\right) \ldots \rho\left(a_{n}\right):=\rho\left(a_{1}\right) \otimes \ldots \otimes \rho\left(a_{n}\right)$.

[^14]:    ${ }^{1}$ Here $\left\langle\left\langle\alpha_{1}, \alpha_{2}, \ldots, \alpha_{n}\right\rangle\right\rangle$ denotes the Pfister form $\left\langle 1, \alpha_{1}\right\rangle \otimes\left\langle 1, \alpha_{2}\right\rangle \otimes \ldots \otimes\left\langle 1, \alpha_{n}\right\rangle$.

[^15]:    ${ }^{2}$ We are doing a convenient use of Corollary 6.2.10

[^16]:    ${ }^{1}$ We are unable to solve this question with the methods so far developed. We believe that to address this question, we will have to develop the theory of quadratic extensions of pre-special hyperfields.

[^17]:    ${ }^{2}$ I.e. for all $\left.a \in G,\langle a, a\rangle \equiv<1,1\right\rangle$ iff $a=1$, but eventually $-1=1$.

