Unification in Commutative Theories, Hilbert's Basis Theorem, and Gröbner Bases

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# Unification in Commutative Theories, Hilbert's Basis Theorem and Gröbner Bases 1) 

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#### Abstract

Unification in a commutative theory E may be reduced to solving linear equations in the corresponding semiring $S(E)$ ( Nutt (1988) ). The unification type of $E$ can thus be characterized by algebraic properties of $\mathrm{S}(\mathrm{E})$. The theory of abelian groups with n commuting homomorphisms corresponds to the semiring $\mathbb{Z}\left[\mathrm{X}_{1}, \ldots, \mathrm{X}_{\mathrm{n}}\right]$. Thus Hilbert's Basis Theorem can be used to show that this theory is unitary. But this argument does not yield a unification algorithm. Linear equations in $\mathbb{Z}\left[\mathrm{X}_{1}, \ldots, \mathrm{X}_{\mathrm{n}}\right]$ can be solved with the help of Gröbner Base methods, which thus provide the desired algorithm. The theory of abelian monoids with a homomorphism is of type zero ( Baader (1988) ). This can also be proved by using the fact that the corresponding semiring, namely $\mathbf{N}[\mathrm{X}]$, is not noetherian. An other example of a semiring (even ring), which is not noetherian, is the ring $\mathbb{Z}<X_{1}, \ldots, X_{n}>$, where $X_{1}, \ldots, X_{n}(n>1)$ are non-commuting indeterminates. This semiring corresponds to the theory of abelian groups with n non-commuting homomorphisms. Surprisingly, by construction of a Gröbner Base algorithm for right ideals in $\mathbb{Z}<X_{1}, \ldots, X_{n}>$, it can be shown that this theory is unitary unifying.


## 1. Introduction

E-unification is concerned with solving term equations modulo an equational theory E . More formally, let E be an equational theory and $=_{E}$ be the equality of terms, induced by E. An $E$-unification problem $\Gamma$ is a finite set of equations $\left\langle s_{i}=t_{i} ; 1 \leq i \leq n\right\rangle_{E}$ where $s_{i}$ and $t_{i}$ are terms. A substitution $\theta$ is called an $E$-unifier of $\Gamma$ iff $s_{i} \theta=E t_{i} \theta$ for each $i, i=1$, $\ldots, n$. The set of all E -unifiers of $\Gamma$ is denoted by $\mathrm{U}_{\mathrm{E}}(\Gamma)$.
In general we do not need the set of all E-unifiers. A complete set of E-unifiers, i.e. a set of E-unifiers from which all E-unifiers may be generated by E-instantiation, is sufficient. More precisely, we extend $=_{E}$ to $U_{E}(\Gamma)$ and define a quasi-ordering $\leq_{E}$ on $U_{E}(\Gamma)$ by
$\sigma=E \theta$ iff $x \sigma=E x \theta$ for all variables $x$ occurring in $s_{i}$ or $t_{i}$ for some $i, i=1, \ldots, n$,
$\sigma \leq_{\mathrm{E}} \theta$ iff there exists a substitution $\lambda$ such that $\sigma=_{\mathrm{E}} \theta \circ \lambda$.
In this case $\sigma$ is called an E-instance of $\theta$.
A complete set $\mathrm{cU}_{\mathrm{E}}(\Gamma)$ of $E$-unifiers of $\Gamma$ is defined as
(1) $\mathrm{cU}_{\mathrm{E}}(\Gamma) \subseteq \mathrm{U}_{\mathrm{E}}(\Gamma)$,
(2) For all $\theta \in U_{E}(\Gamma)$ there exists $\sigma \in \mathrm{cU}_{\mathrm{E}}(\Gamma)$ such that $\theta \leq_{E} \sigma$.

[^0]For reasons of efficiency this set should be as small as possible. Thus we are interested in minimal complete sets of E-unifiers, that means complete sets where two different elements are not comparable w.r.t. E-instantiation. The unification type of a theory E is defined with reference to the cardinality and existence of minimal complete sets. The theory E is unitary ( finitary, infinitary ) iff minimal complete sets of E -unifiers always exist and their cardinality is at most one ( always finite, at least once infinite). E has unification type zero iff there is an E-unification problem without minimal complete set of E-unifiers. If the terms may contain free constants, we talk about unification with constants, else about unification without constants ( see Baader (1988), Section 7 ). If nothing else is specified, "unification" means "unification without constants". For more information about unification theory and the unification hierarchy consult Siekmann (1988).
Unification in the empty theory ( which is unitary ) plays an important rôle in automated theorem proving, term rewriting and logic programming. Generalizations to E-unification usually require that E is finitary ( see e.g. Stickel (1985), Jouannaud-Kirchner (1986) and Jaffar-Lassez-Maher (1984) ). A finitary theory most used in this context is the theory of abelian semigroups ( monoids), i.e. the theory of an associative, commutative binary operation ( with a neutral element ). Unification algorithms for this theory ( see e.g. Live-sey-Siekmann (1978), Stickel (1981), Fages (1984), Fortenbacher (1985), Büttner (1986), Herold (1987) ) make use of the fact that unifiers correspond to solutions of systems of linear equations in the semiring $\mathbb{N}$ ( see Eilenberg (1974) or Kuich-Salomaa (1986) for the definition and properties of semirings ). The same phenomenon occurs for the theory of abelian groups where the semiring is $\mathbb{Z}$ ( Lankford-Butler-Brady (1984)) and for the theory of idempotent abelian monoids where the 2-element boolean semiring $\mathcal{B}$ is used (Livesey-Siekmann (1978), Baader-Büttner (1988) ).
These three theories belong to the class of commutative theories ( roughly speaking, theories where the finitely generated free objects are direct products of the free objects in one generator ), which were defined in Baader (1988). In that paper it is shown that con-stant-free unification in commutative theories is either unitary or of type zero and there are given sufficient conditions for a commutative theory to be unitary ( resp. finitary w.r.t. unification with constants ). The above mentioned results for abelian monoids etc. and some new results ( for abelian monoids with an involution, idempotent abelian monoids with an involution, abelian groups with an involution, abelian groups of exponent m ) could thus be obtained as corollaries to a general theorem. In Baader (1989) these conditions were modified to algebraic characterizations of unification type unitary for constantfree unification and type finitary for unification with constants in commutative theories. An interesting consequence of these characterizations is the fact that commutative theories are always unitary ( finitary w.r.t. unification with constants ), if the finitely generated free objects are finite (Baader (1988)).
Werner Nutt ( Nutt (1988) ) observed that commutative theories are (modulo a translation of the signature ) what he calls monoidal theories and that unification in these theories may always be reduced to solving linear equations in certain semirings. He pointed out that the theory of abelian groups with a homomorphism corresponds to the semiring $\mathbb{Z}[\mathrm{X}]$. Thus Hilbert's Basis Theorem can be used to prove that the theory of abelian groups with a homomorphism is unitary. But this argument does not yield a unification algorithm. Linear equations in $\mathbb{Z}[\mathrm{X}]$ can be solved with the help of Gröbner Base methods ( see Buchberger (1985) and Section 6 of this paper ), which thus provide the desired algorithm.
The theory of abelian monoids with a homomorphism is of type zero ( Baader (1988) ). This can also be demonstrated using the fact that the corresponding semiring, namely
$\operatorname{IN}[\mathrm{X}]$, is not noetherian (Section 4 ).
Another example of a semiring which is not noetherian is the ring $\mathbb{Z}\langle X, Y\rangle$, where $\mathrm{X}, \mathrm{Y}$ are non-commuting indeterminates. This semiring corresponds to the theory of abelian groups with two ( non-commuting ) homomorphisms. Surprisingly, by construction of a Gröbner Base algorithm for right ideals in $\mathbb{Z}<X, Y>$, I was able to show that this theory is unitary unifying. Of course, this result can be extended to an arbitrary, finite number of non-commuting indeterminates (Section 8 and 9 ).

## 2. Commutative Theories

In this section we give a definition of commutative theories, recall some of the properties derived in Baader (1988) and show how the corresponding semirings may be obtained in this framework.
An equational theory E defines a variety $V(E)$, i.e. the class of all algebras (of the given signature $\Omega$ ) which satisfy each identity of $E$. For any set $X$ of generators, $V(E)$ contains a free algebra over $V(E)$ with generators $X$, which will be denoted by $\mathrm{F}_{\mathrm{E}}(\mathrm{X})$.
Let $F(E)$ be the class of all free algebras $F_{E}(X)$ with finite sets $X$ and let $C(E)$ be the category which has the elements of $F(E)$ as objects and the homomorphisms between these elements as morphisms. Note that the coproduct of $\mathrm{F}_{\mathrm{E}}(\mathrm{X})$ and $\mathrm{F}_{\mathrm{E}}(\mathrm{Y})$ in $\mathrm{C}(\mathrm{E})$ is given by $\mathrm{F}_{\mathrm{E}}(\mathrm{X} \bullet \mathrm{\cup})$ ( where $\bullet^{\bullet}$ means disjoint union ). Thus $\mathrm{F}_{\mathrm{E}}(\mathrm{X})$ is the coproduct of the isomorphic objects $\mathrm{F}_{\mathrm{E}}(\mathrm{x})$ for $\mathrm{x} \in \mathrm{X}$.

Let $\Gamma=\left\langle\mathrm{s}_{\mathrm{i}}=\mathrm{t}_{\mathrm{i}} ; 1 \leq \mathrm{i} \leq \mathrm{n}\right\rangle_{\mathrm{E}}$ be an E-unification problem and X be the (finite) set of variables $x$ occurring in some $s_{i}$ or $t_{i}$. Evidently we can consider the $s_{i}$ and $t_{i}$ as elements of $F_{E}(X)$. Since we do not distinguish between $=_{E}$-equivalent unifiers, any $E$-unifier of $\Gamma$ can be regarded as a homomorphism of $\mathrm{F}_{\mathrm{E}}(\mathrm{X})$ into $\mathrm{F}_{\mathrm{E}}(\mathrm{Y})$ for some finite set Y ( of variables ). Let $\mathrm{I}=\left\{\mathrm{x}_{1}, \ldots, \mathrm{x}_{\mathrm{n}}\right\}$ be a set of cardinality n . We define homomorphisms

$$
\sigma, \tau: \mathrm{F}_{\mathrm{E}}(\mathrm{I}) \rightarrow \mathrm{F}_{\mathrm{E}}(\mathrm{X}) \text { by } \mathrm{x}_{\mathrm{i}} \sigma:=\mathrm{s}_{\mathrm{i}} \text { and } \mathrm{x}_{\mathrm{i}} \tau:=\mathrm{t}_{\mathrm{i}} \quad(\mathrm{i}=1, \ldots, \mathrm{n}) .
$$

Now $\delta: F_{E}(X) \rightarrow F_{E}(Y)$ is an E-unifier of $\Gamma$ iff $x_{i} \sigma \delta=s_{i} \delta=t_{i} \delta=x_{i} \tau \delta$ for $i=1, \ldots$, n, i.e. iff $\sigma \delta=\tau \delta$. Thus an E-unification problem can be written as a pair $\langle\sigma=\tau\rangle_{\mathrm{E}}$ of morphisms $\sigma, \tau: \mathrm{F}_{\mathrm{E}}(\mathrm{I}) \rightarrow \mathrm{F}_{\mathrm{E}}(\mathrm{X})$ in the category $\mathrm{C}(\mathrm{E})$. An E-unifiers of the unification problem $\langle\sigma=\tau\rangle_{\mathrm{E}}$ is a morphism $\delta$ such that $\sigma \delta=\tau \delta$.
This categorical reformulation of E-unification ( due to Rydeheard-Burstall (1985) ) allows to characterize the class of commutative theories by properties of the category $\mathrm{C}(\mathrm{E})$ of finitely generated E -free objects: $\mathrm{C}(\mathrm{E})$ has to be a semiadditive category ( see Herrlich-Strecker (1973) and Baader (1988) ). In order to give a more algebraic definition of commutative theories we need some more notation.
A constant symbol (i.e. a nullary function symbol) e $\in \Omega$ is called idempotent in $E$ iff for any $f \in \Omega$ we have $f(e, \ldots, e)=E e$, i.e. in any algebra $A \in V(E), f(e, \ldots, e)=e$ holds. Note that for nullary $f$ this means $f=E$.
Let $\mathbf{K}$ be a class of algebras ( of signature $\Omega$ ). An $n$-ary implicit operation in $\mathbf{K}$ is a family $f=\left\{f_{A} ; A \in K\right\}$ of mappings $f_{A}: A^{n} \rightarrow A$ which is compatible with all homomor-
phisms, i.e. for any homomorphism $h: A \rightarrow B$ with $A, B \in K$ and all $a_{1}, \ldots, a_{n} \in A$, $f_{A}\left(a_{1}, \ldots, a_{n}\right) h=f_{B}\left(a_{1} h, \ldots, a_{n} h\right)$ holds. In the following we omit the index and just write $f$ for any $f_{A}$. Obviously an $\Omega$-term induces an implicit operation on any class of $\Omega$-algebras.

DEFINITION 2.1. An equational theory E is called commutative iff the following holds:
(1) $\Omega$ contains a constant symbol $e$, which is idempotent in $E$.
(2) There is a binary implicit operation $*$ in $\mathrm{F}(\mathrm{E})$ such that
(a) The constant $e$ is a neutral element for $*$ in any algebra $A \in F(E)$.
(b) For any $n$-ary function symbol $f \in \Omega$, any algebra $A \in F(E)$ and any $s_{1}, \ldots, s_{n}, t_{1}$, $\ldots, \mathrm{t}_{\mathrm{n}} \in \mathrm{A}$ we have $\mathrm{f}\left(\mathrm{s}_{1} * \mathrm{t}_{1}, \ldots, \mathrm{~s}_{\mathrm{n}} * \mathrm{t}_{\mathrm{n}}\right)=\mathrm{f}\left(\mathrm{s}_{1}, \ldots, \mathrm{~s}_{\mathrm{n}}\right) * \mathrm{f}\left(\mathrm{t}_{1}, \ldots, \mathrm{t}_{\mathrm{n}}\right)$.

In Baader (1988) the following properties for commutative theories E are shown within a categorical framework, using well-known results for semiadditive categories:
(2.2) $\left|F_{E}(\varnothing)\right|=1$ and $F_{E}(\varnothing)$ is the zero object of $C(E)$.
(2.3) The implicit operation $*$ of Definition 2.1 is associative and commutative. It induces a binary operation + on any morphism set $\operatorname{hom}\left(\mathrm{F}_{\mathrm{E}}(\mathrm{X}), \mathrm{F}_{\mathrm{E}}(\mathrm{Y})\right)$ as follows: Let $\sigma, \tau$ : $\mathrm{F}_{\mathrm{E}}(\mathrm{X}) \rightarrow \mathrm{F}_{\mathrm{E}}(\mathrm{Y})$ and $\mathrm{s} \in \mathrm{F}_{\mathrm{E}}(\mathrm{X})$. Then $\mathrm{s}(\sigma+\tau):=(\mathrm{s} \sigma) *(\mathrm{~s} \tau)$.
This operation is also associative and commutative and it distributes with the composition of morphisms. The morphism $0: F_{E}(X) \rightarrow F_{E}(Y)$ defined by $x \mapsto e$ for all $x$ $\in X$ is the zero morphism in hom $\left(F_{E}(X), F_{E}(Y)\right)$ and it is a neutral element for + on $\operatorname{hom}\left(\mathrm{F}_{\mathrm{E}}(\mathrm{X}), \mathrm{F}_{\mathrm{E}}(\mathrm{Y})\right.$ )
(2.4) The coproduct $\mathrm{F}_{\mathrm{E}}(\mathrm{X} \cup \mathrm{Y})$ of $\mathrm{F}_{\mathrm{E}}(\mathrm{X})$ and $\mathrm{F}_{\mathrm{E}}(\mathrm{Y})$ is also the product of these objects, i.e. $F_{E}(X \cup Y) \cong F_{E}(X) \times F_{E}(Y)$.
(2.5) Consider $\sigma: F_{E}(X) \rightarrow F_{E}(Y)$. Let $u_{x}$ for $x \in X\left(p_{y}\right.$ for $\left.y \in Y\right)$ be the injections of the coproduct $\mathrm{F}_{\mathrm{E}}(\mathrm{X})$ ( projections of the product $\mathrm{F}_{\mathrm{E}}(\mathrm{Y})$ ). Then $\sigma$ is uniquely determined by the matrix $M_{\sigma}=\left(u_{x} \sigma p_{y}\right)_{x \in X, y \in Y}$. For $\sigma, \tau: F_{E}(X) \rightarrow F_{E}(Y)$ and $\delta$ : $\mathrm{F}_{\mathrm{E}}(\mathrm{Y}) \rightarrow \mathrm{F}_{\mathrm{E}}(\mathrm{Z})$ we have $\mathrm{M}_{\sigma+\tau}=\mathrm{M}_{\sigma}+\mathrm{M}_{\tau}$ and $\mathrm{M}_{\sigma \delta}=\mathrm{M}_{\sigma} \cdot \mathrm{M}_{\delta}$.

Werner Nutt ( Nutt (1988) ) observed that commutative theories are (modulo a translation of the signature ) what he calls monoidal theories and that unification in a monoidal theory E may be reduced to solving linear equations in a certain semiring $\mathrm{S}(\mathrm{E})$. In our framework this semiring can be obtained as follows:
Let 1 be an arbitrary set of cardinality 1. Property (2.3) yields that hom $\left(\mathrm{F}_{\mathrm{E}}(\mathbf{1}), \mathrm{F}_{\mathrm{E}}(\mathbf{1})\right.$ ) with addition " + " and composition as multiplication is a semiring, which shall be denoted by $S(E)$. Any $\mathrm{F}_{\mathrm{E}}(\mathrm{x})$ is isomorphic to $\mathrm{F}_{\mathrm{E}}(1)$ and for $|\mathrm{X}|=\mathrm{n}, \mathrm{F}_{\mathrm{E}}(\mathrm{X})$ is n-th power and copower of $\mathrm{F}_{\mathrm{E}}(\mathbf{1})$. Thus, for $\sigma: \mathrm{F}_{\mathrm{E}}(\mathrm{X}) \rightarrow \mathrm{F}_{\mathrm{E}}(\mathrm{Y})$, the entries $u_{\mathrm{x}} \sigma \mathrm{p}_{\mathrm{y}}$ of the $|\mathrm{X}| \times|Y|$-matrix $M_{\sigma}$ may all be considered as elements of $S(E)$. Hence all morphisms of $C(E)$ can be written as matrices over the semiring $\mathrm{S}(\mathrm{E})$. Addition and composition of morphisms correspond to addition and multiplication of matrices over $S(E)$ as stated in (2.5).

We now give some examples of commutative theories, whose unification properties will be considered in subsequent sections of this paper. In all these examples, the implicit operation is given by a function symbol, which is associative and commutative in the corresponding theory. Additional examples of commutative theories can be found in Baader (1988).

EXAMPLES 2.6. We consider the following signatures:
$\Sigma:=\{\cdot, 1, h\}$, where $\cdot$ is binary, 1 is nullary and $h$ is unary.
For $n \geq 0, \Omega_{n}:=\left\{\cdot, 1,{ }^{-1}, h_{1}, \ldots, h_{n}\right\}$, where $\cdot$ is binary, 1 is nullary and ${ }^{-1}$ and the $h_{i}$ are unary.
(1) The theory AMH of abelian monoids with a homomorphism. The signature is $\Sigma$ and $\mathrm{AMH}:=\{x \cdot 1=x, x \cdot(y \cdot z)=(x \cdot y) \cdot z, x \cdot y=y \cdot x$,

$$
h(x \cdot y)=h(x) \cdot h(y), h(1)=1\}
$$

(2) The theory AIMH of idempotent abelian monoids with a homomorphism. The signature is $\Sigma$ and AIMH $:=$ AMH $\cup\{x \cdot x=x\}$.
(3) The theory AGnH of abelian groups with $n$ (non-commuting) homomorphisms. We take signature $\Omega_{\mathrm{n}}$ and define $\mathrm{AGnH}:=\{\mathrm{x} \cdot 1=\mathrm{x}, \mathrm{x} \cdot(\mathrm{y} \cdot \mathrm{z})=(\mathrm{x} \cdot \mathrm{y}) \cdot \mathrm{z}, \mathrm{x} \cdot \mathrm{y}=\mathrm{y} \cdot \mathrm{x}$, $\left.\mathrm{x} \cdot \mathrm{x}^{-1}=1\right\} \cup\left\{\mathrm{h}_{\mathrm{i}}(\mathrm{x} \cdot \mathrm{y})=\mathrm{h}_{\mathrm{i}}(\mathrm{x}) \cdot \mathrm{h}_{\mathrm{i}}(\mathrm{y}) ; 1 \leq \mathrm{i} \leq \mathrm{n}\right\}$.
(4) The theory AGnHC of abelian groups with $n$ commuting homomorphisms. The signature is $\Omega_{\mathrm{n}}$ and AGnHC $:=\mathrm{AGnH} \cup\left\{\mathrm{h}_{\mathrm{i}}\left(\mathrm{h}_{\mathrm{j}}(\mathrm{x})\right)=\mathrm{h}_{\mathrm{j}}\left(\mathrm{h}_{\mathrm{i}}(\mathrm{x})\right) ; 1 \leq \mathrm{i}<\mathrm{j} \leq \mathrm{n}\right\}$.

It is easy to see that these theories are commutative. Note that the implicit operation induced by the term $x \cdot y$ (for a binary function symbol $\cdot$ ) satisfies $2 b$ of Definition 2.1 for $f=\cdot \operatorname{iff}(a \cdot b) \cdot(c \cdot d)=(a \cdot c) \cdot(b \cdot d)$ holds in any algebra $A \in F(E)$ and for $f=h($ for $a$ unary function symbol $h$ ) iff $h(x \cdot y)=h(x) \cdot h(y)$ holds.

## 3. Unification in Commutative Theories

In this section we recall the characterizations of unification type unitary (finitary for unification with constants ) for commutative theories given in Baader (1989). As a consequence we derive, that unification in a commutative theory E means solving systems of linear equations in the semiring $S(E)$. This yields algebraic characterizations of the unification types.

THEOREM 3.1. A commutative theory E is unitary iff it satisfies the following condition:
Let $y$ be an arbitrary variable. For any E-unification problem $\langle\sigma=\tau\rangle_{\mathrm{E}}$ ( where $\sigma, \tau$ : $\left.F_{E}(I) \rightarrow F_{E}(X)\right)$ there are finitely many E-unifiers $\alpha_{1}, \ldots, \alpha_{r}: F_{E}(X) \rightarrow F_{E}(y)$ such that any E-unifier $\delta: F_{E}(X) \rightarrow F_{E}(y)$ is representable as

$$
\delta=\sum_{i=1}^{i=r} \alpha_{i} \lambda_{i}
$$

where $\lambda_{i}: F_{E}(y) \rightarrow F_{E}(y)$ are morphisms.
If we translate morphisms into matrices over $\mathrm{S}(\mathrm{E})$, we obtain the following reformulation of Theorem 3.1:

COROLLARY 3.2. A commutative theory E is unitary iff the corre $\mathrm{c}_{\mathrm{I}}$ onding semiring $S(E)$ satisfies the following condition: For any $n, m \geq 1$ and any pair $M_{\sigma}, M_{\tau}$ of $m \times n$-matrices over $\mathrm{S}(\mathrm{E})$ the set

$$
\mathrm{U}\left(\mathrm{M}_{\sigma}, \mathrm{M}_{\tau}\right):=\left\{\underline{x} \in \mathrm{~S}(E)^{\mathrm{n}} ; \mathrm{M}_{\sigma \underline{x}}=\mathrm{M}_{\tau} \underline{x}\right\}
$$

is a finitely generated right $S(E)$-semimodule, i.e. there are finitely many $\underline{x}_{1}, \ldots, \underline{x}_{\mathrm{r}} \in$ $S(E)^{n}$ such that $U\left(M_{\sigma}, M_{\tau}\right)=\left\{\underline{x}_{1} s_{1}+\ldots+\underline{x}_{\uparrow} s_{r} ; s_{1}, \ldots, s_{r} \in S(E)\right\}$.

THEOREM 3.3. Let E be a unitary commutative theory. Then E is finitary w.r.t. unification with constants iff the following condition holds:

For any morphism ( of $C(E)$ ) $\delta: F_{E}(X) \rightarrow F_{E}(Y)$ there exist finite sets $M, N$ such that:
(1) The elements of $M$ are morphisms $\mu: F_{E}(Y) \rightarrow F_{E}(X)$ satisfying $\delta \mu=1$.
(2) The elements of $N=\left\{v_{1}, \ldots, v_{r}\right\}$ are morphisms $v_{i}: F_{E}(Y) \rightarrow F_{E}\left(Z_{i}\right)$ with $\delta v_{i}=0$.
(3) For any $\lambda: \mathrm{F}_{\mathrm{E}}(\mathrm{Y}) \rightarrow \mathrm{F}_{\mathrm{E}}(\mathrm{X})$ with $\delta \lambda=1$ there are $\mu \in \mathrm{M}$ and morphisms $\lambda_{1}, \ldots, \lambda_{\mathrm{T}}$ ( where $\lambda_{i}: \mathrm{F}_{\mathrm{E}}\left(\mathrm{Z}_{\mathrm{i}}\right) \rightarrow \mathrm{F}_{\mathrm{E}}(\mathrm{X})$ ) satisfying

$$
\lambda=\mu+\sum_{\mathrm{i}=1}^{\mathrm{i}=\mathrm{r}} v_{\mathrm{i}} \lambda_{\mathrm{i}}
$$

The translation of morphisms into matrices over $S(E)$ yields a sufficient condition for $E$ to be finitary w.r.t. unification with constants.

COROLLARY 3.4. Let E be a unitary commutative theory. Then E is finitary w.r.t. unification with constants, if the following condition holds in $S(E)$ :
Let $A$ be any $m \times n$-matrices over $S(E)$ and let $\underline{b}$ be any element of $S(E)^{m}$. Then the set $M:=\left\{\underline{x} \in S(E)^{n} ; A \underline{x}=\underline{b}\right\}$ is a finite union of cosets of the (finitely generated) right $S(E)$-semimodule $N:=\left\{\underline{x} \in S(E)^{n} ; A \underline{x}=0\right\}$, i.e. there exist finitely many $\underline{m}_{1}, \ldots, \underline{m}_{k} \in$ $\mathrm{S}(\mathrm{E})^{\mathrm{n}}$ such that $\mathrm{M}=\left\{\underline{m}_{\mathrm{i}}+\underline{\mathrm{n}} ; \underline{\mathrm{n}} \in \mathrm{N}\right.$ and $\left.1 \leq \mathrm{i} \leq \mathrm{k}\right\}$.

Note that the semimodule $N$ is finitely generated, since $E$ is unitary and $N=U(A, 0)$, where 0 is the $m \times n$ zero matrix. From Theorem 3.3 we can only deduce, that the condition of the corollary is sufficient, since in Theorem 3.3 we talk about specific inhomogeneous equations $\mathrm{AX}=\mathrm{E}$, while in Corollary 3.4 the right-hand side of the equation is an arbitrary vector b .
Assume that $\mathrm{S}(\mathrm{E})$ is a ring and let $\underline{x}_{0}$ be an arbitrary solution of the inhomogeneous equation $A \underline{x}=\underline{b}$. Then any solution $y$ of $A \underline{x}=\underline{b}$ is of the form $y=\underline{x}_{0}+\underline{z}$, where $\underline{z}:=y-\underline{x}_{0}$ is a solution of the homogeneous equation $\mathrm{Ax}=0$. This proves

COROLLARY 3.5. Let $E$ be a unitary commutative theory such that $S(E)$ is a ring. Then E is unitary w.r.t. unification with constants.

## 4. A Commutative Theory of Unification Type Zero

In 1972 Plotkin conjectured, that there exists an equational theory $E$ which has unification type zero. But only in 1983, Fages and Huet constructed the first example of an equational theory of this type. Schmidt-Schauß (1986) and the present author (1986) showed that the theory of idempotent semigroups is of unification type zero and in Baader (1987) I have proved, that almost all varieties of idempotent semigroups are defined by type zero theories. This provides us with countably many examples of type zero theories, which are more natural than the original example of Fages and Huet.
In Baader (1988) it is shown that the theory AIMH of idempotent abelian monoids with a homomorphism has type zero. The same proof can be used for AMH, the theory of abelian monoids with a homomorphism, in place of AIMH. This section contains a more algebraic proof of the fact that AMH has type zero. Since commutative theories are either unitary or of unification type zero (Baader (1988), Theorem 6.1 ), it is sufficient to show, that the semiring $\mathrm{S}(\mathrm{AMH})$ does not satisfy the condition of Corollary 3.2.
Let $\sigma: \mathrm{F}_{\mathrm{AMH}}(\mathrm{x}) \rightarrow \mathrm{F}_{\mathrm{AMH}}(\mathrm{x})$ be a morphism of $\mathrm{C}(\mathrm{AMH})$. Then there are $\mathrm{k} \geq 0$ and $\mathrm{a}_{0}$, $\ldots, a_{k} \in \mathbb{N}$ such that $x \sigma=x^{a_{0}} h\left(x^{a_{1}}\right) \ldots h^{k}\left(x^{a_{k}}\right)$. We associate with the morphism $\sigma$ the polynomial $p_{\sigma}=a_{0}+a_{1} X+\ldots+a_{k} X^{k} \in \mathbb{N}[X]$. It is easy to see that $p_{\sigma \delta}=p_{\sigma} \cdot p_{\delta}$ and $p_{\sigma+\delta}=$ $p_{\sigma}+p_{\delta}$, which shows that $S(A M H) \equiv \mathbb{N}[X]$.
We consider the linear equation (*) $X x_{1}+X x_{2}=x_{2}+X^{2} x_{3}$, which has to be solved by a vector $\underline{p}=\left(p_{1}, p_{2}, p_{3}\right)$ in $(\mathbb{N}[X])^{3}$. Obviously, for any $n \geq 0$, the vector $p^{(n)}=\left(p_{1}^{(n)}, p_{2}^{(n)}\right.$, $\left.p_{3}^{(n)}\right)=\left(1, X+X^{2}+\ldots+X^{n+1}, X^{n}\right)$ is a solution of $(*)$.

LEMMA 4.1. There does not exist a solution $p$ of $(*)$ in $(\mathbb{N}[\mathrm{X}])^{3}$ such that $\mathrm{p}_{1}+\mathrm{p}_{3}=1$.
PROOF. For $p_{1}=0$ and $p_{3}=1$ we get $X p_{2}=p_{2}+X^{2}$, which yields $(X-1) p_{2}=X^{2}$ in $\mathbb{Z}[X]$. But $X-1$ is not a divisor of $X^{2}$. The case $p_{1}=1$ and $p_{3}=0$ leads to a similar contradiction.

It is easy to see that $I_{1+3}:=\left\{p_{1}+p_{3}\right.$; There exists $p_{2}$ such that $\left(p_{1}, p_{2}, p_{3}\right)$ solves (*) \} is an ideal in $\mathbb{N}[X]$. We know that $1+\mathrm{X}^{\mathrm{n}} \in \mathrm{I}_{1+3}$ for any $\mathrm{n} \geq 0$ and $1 \notin \mathrm{I}_{1+3}$.

## LEMMA 4.2.

An ideal $I \subseteq \mathbb{N}[X]$ such that $1+X^{n} \in I$ for any $n \geq 0$ and $1 \notin I$ is not finitely generated.
PROOF. Evidently $1+X^{n}=f \cdot g$ for $f, g \in \mathbb{N}[X]$ or $1+X^{n}=f+g$ for $f, g \in \mathbb{N}[X] \backslash\{0\}$ implies $f=1$ or $g=1$. But $1 \notin I$.

PROPOSITION 4.3. The theory AMH has unification type zero.
PROOF. Assume that AMH has not type zero. Then AMH is unitary and, by Corollary $3.2, \underline{I}:=\left\{\underline{p} \in(\mathbb{N}[X])^{3} ; \underline{p}\right.$ is a solution of $\left.(*)\right\}$ is a finitely generated right $\mathbb{N}[X]$-semimodule. But then $\mathrm{I}_{1+3}=\left\{\mathrm{p}_{1}+\mathrm{p}_{3}\right.$; There exists $\mathrm{p}_{2}$ such that $\left.\left(\mathrm{p}_{1}, \mathrm{p}_{2}, \mathrm{p}_{3}\right) \in \underline{I}\right\}$ would also be finitely generated, which contradicts Lemma 4.2.

The fact that the set of solutions of the equation (*) is not a finitely generated right semi-
module is not specific for the semiring $\mathbb{N}[\mathrm{X}]$. More general, let $S$ be a semiring which is not a ring ( that means, that there exists $s \in S$ such that for all $t \in S s+t \neq 0$ ). Then the right $S[X]$-semimodule $I:=\left\{p \in(S[X])^{3} ; p\right.$ is a solution of $\left.(*)\right\}$ is not finitely generated (Baader-Nutt (1989) ).

## 5. AGnHC-Unification and Hilbert's Basis Theorem

It is easy to see that $\mathrm{S}(\mathrm{AGnHC})$ is isomorphic to the ring $\mathbb{Z}\left[X_{1}, \ldots, X_{n}\right]$, i.e. the polynomial ring over $\mathbb{Z}$ in the (commuting ) indeterminates $X_{1}, \ldots, X_{n}$. To establish the condition of Corollary 3.2, we have to consider systems of homogeneous linear equations in $\mathbb{Z}$ [ $X_{1}$, $\left.\ldots, X_{n}\right]$, i.e. systems $f_{1 i} x_{1}+\ldots+f_{k i} x_{k}=0(i=1, \ldots, s)$, where the coefficients $f_{i j}$ and the desired solutions are elements of $\mathbb{Z}\left[X_{1}, \ldots, X_{n}\right]$. The set of solutions $\underline{I} \subseteq\left(\mathbb{Z}\left[X_{1}, \ldots, X_{n}\right]\right)^{k}$ is a $\mathbb{Z}\left[X_{1}, \ldots, X_{n}\right]$-module, which is finitely generated by Hilbert's Basis Theorem and the fact that $\mathbb{Z}$ is a noetherian ring ( see e.g. Jacobson (1980)). Thüs AGnHC is unitary w.r.t. unification without constants. Since $\mathbb{Z}\left[X_{1}, \ldots, X_{n}\right]$ is a ring, Corollary 3.5 applies and we have proved

PROPOSITION 5.1. ( Nutt (1988))
For any $\mathrm{n} \geq 0$ the theory AGnHC is unitary and it is also unitary w.r.t. unification with constants.

This proof of Proposition 5.1 does not yield an AGnHC-unification algorithm, because we still do not know how to solve linear equations in $\mathbb{Z}\left[\mathrm{X}_{1}, \ldots, \mathrm{X}_{\mathrm{n}}\right]$ effectively. The next section describes one possible answer to this problem.

## 6. Solving Linear Equations in $\mathbb{Z}\left[X_{1}, \ldots, X_{n}\right]$ using Gröbner Bases

Buchberger (1985) describes an effective method, which constructs finitely many generators of the solutions of a single equation $f_{1} x_{1}+\ldots+f_{k} x_{k}=0$, where the $f_{i}$ and the desired solutions are elements of $K\left[X_{1}, \ldots, X_{n}\right]$ for a field $K$. This method may also be used for $\mathbb{Z}\left[X_{1}, \ldots, X_{n}\right]$ ( see Buchberger (1985) for Gröbner Bases of polynomials over $\mathbb{Z}$ and Kandry-Rody-Kapur (1988) for Gröbner Bases of polynomials over a euclidean ring ), but the proof of its correctness becomes more involved. Systems of equations can then be solved by successive substitution. A more efficient approach to solving systems of equations is described in Furukawa-Sasaki-Kobayashi (1986), where Gröbner base theory is extended to modules over $K\left[X_{1}, \ldots, X_{n}\right]$.

First we recall some facts and notations concerning Gröbner bases:
(6.1) Admissible term orderings.

Let $T_{n}:=\left\{X_{1}^{k_{1}} \ldots X_{n}^{k_{n}} ; k_{1}, \ldots, k_{n} \in \mathbb{N}\right\}$ be the set of all terms (i.e. monomials with coefficient 1 ) in $\mathbb{Z}\left[X_{1}, \ldots, X_{n}\right]$. With respect to multiplication of polynomials, $T_{n}$ is a commuta-
tive monoid ( with neutral element $1=X_{1}^{0} \ldots X_{n}^{0}$ ), which is isomorpic to the additive monoid $\mathbb{N}^{\mathrm{n}}$.
A linear ordering $<$ on $T_{n}$ is called compatible iff for all $r, s, t \in T_{n} r<s$ implies $r t<s t$ and it is called admissible iff it is compatible and satisfies $1<\mathrm{s}$ for all $\mathrm{s} \in \mathrm{T}_{\mathrm{n}}$. It is easy to see that a compatible linear ordering on $T_{n}$ is admissible iff it is noetherian.
Complete descriptions of all compatible linear orderings have been given by Trevisan (1953), Zaiceva (1953) and, more recently, by Robbiano (1985) and Martin (1988):

Any compatible linear ordering < on $T_{n}$ is completely determined by a $n \times s$ matrix $U_{<}$of $s$ $\leq n$ orthogonal vectors $u_{1}, \ldots, u_{s} \in \mathbb{R}^{n}$ of $\mathbb{Q}$-dimension $n$ as follows: $X_{1}^{k_{1}} \ldots X_{n}^{k_{n}}<X_{1}^{h_{1}} \ldots X_{n}^{h_{n}}$ iff the first non-zero element of $\left(h_{1}-k_{1}, \ldots, h_{n}-k_{n}\right) \cdot U_{<}$is greater than zero.
It is easy to see that the compatible linear ordering < is admissible iff in any row of $\mathrm{U}_{<}$, the first non-zero entry is greater than zero.

## (6.2) Rewriting with polynomials.

For a polynomial $f$ and a term $t$ which occurs in $f$, coeff $(t, f)$ denotes the coefficient of $t$ in $f$. If $t$ does not occur in $f$, we define coeff( $t, f):=0$. Let $<$ be an admissible ordering and let $f=$ $\mathrm{a} \cdot \mathrm{t}+\mathrm{g}$ be a polynomial in $\mathbb{Z}\left[\mathrm{X}_{1}, \ldots, \mathrm{X}_{\mathrm{n}}\right]$ such that $\mathrm{t} \in \mathrm{T}_{\mathrm{n}}$ is the greatest term in f w.r.t. $<$ and coeff $(t, f)=a \in \mathbb{Z}$ is the coefficient of $t$ in $f$. Then $t$ is called head-term of $f(H T(f))$, a is called head-coefficient of $\mathrm{f}(\mathrm{HC}(\mathrm{f})$ ), a.t is called head-monomial of $\mathrm{f}(\mathrm{HM}(\mathrm{f})$ ) and g $=f-H M(f)$ is called rest of $f(R(f))$.
A set $F$ of polynomials induces the following rewrite relation on $\mathbb{Z}\left[X_{1}, \ldots, X_{n}\right]$ :

## $f \rightarrow_{F} g$ iff (1) f contains a term $t$ with coefficient $a$.

(2) $F$ contains a polynomial $h$ such that $H T(h)=t \cdot s$ ( for some $s \in T_{n}$ ) and $|\mathrm{HC}(\mathrm{h})| \leq \mathrm{lal}$.
(3) $\mathrm{g}=\mathrm{f}-\mathrm{h} \cdot \mathrm{b} \cdot \mathrm{s}$, where $\mathrm{a}=\mathrm{b} \cdot \mathrm{HC}(\mathrm{h})+\mathrm{c}$ for $0 \leq \mathrm{c}<|\mathrm{HC}(\mathrm{h})|, \mathrm{b}, \mathrm{c} \in \mathbb{Z}$.

Let $\stackrel{*}{F}_{F}\left(\right.$ resp. $\dagger_{\mathrm{F}}$ ) denote the reflexive, transitive ( resp. transitive) closuie of $\rightarrow_{\mathrm{F}}$. It can be shown (using a multiset extension of $<$ ) that ${ }^{{ }^{\prime}}{ }_{F}$ is noetherian. The set F generates an ideal $\langle\mathrm{F}\rangle$ in $\mathbb{Z}\left[\mathrm{X}_{1}, \ldots, \mathrm{X}_{\mathrm{n}}\right]$ and this ideal induces a congruence $\equiv_{<\mathrm{F}\rangle}$, namely f $\equiv_{<\mathrm{F}\rangle} \mathrm{g}$ iff $\mathrm{f}-\mathrm{g} \in\langle\mathrm{F}\rangle$. This congruence is the reflexive, transitive and symmetric closure of $\rightarrow_{F}$ (Bachmair-Buchberger (1980)).
(6.3) Gröbner bases and S-polynomials.

Let $I$ be an ideal in $\mathbb{Z}\left[\mathrm{X}_{1}, \ldots, \mathrm{X}_{\mathrm{n}}\right]$ and B let be a finite set of polynomials. B is a Gröbner base for I iff $<\mathrm{B}\rangle=\mathrm{I}$ and $\rightarrow_{\mathrm{B}}$ is confluent. Since ${ }^{{ }^{\prime}}{ }_{\mathrm{B}}$ is noetherian, confluence is equivalent to local confluence and this property can be tested with the help of finitely many critical pairs, which are here called S-polynomials.
Let $g_{1}=c_{1} \cdot t_{1}+R\left(g_{1}\right)$ and $g_{2}=c_{2} \cdot t_{2}+R\left(g_{2}\right)$ be elements of $B$ such that $c_{1} \geq c_{2} \geq 0$ (without loss of generality we assume, that the head coefficients of the polynomials in B are positive $)$. The $S$-polynomial $S\left(g_{1}, g_{2}\right)$ of $g_{1}$ and $g_{2}$ is defined as follows:

Let $\mathrm{s}_{1} \cdot \mathrm{t}_{1}=\mathrm{s}_{2} \cdot \mathrm{t}_{2}=\operatorname{lcm}\left(\mathrm{t}_{1}, \mathrm{t}_{2}\right)$ and $\mathrm{c}_{1}=\mathrm{a} \cdot \mathrm{c}_{2}+\mathrm{b}, 0 \leq \mathrm{b}<\mathrm{c}_{2} \leq \mathrm{c}_{1}, \mathrm{a} \geq 1$. Then

$$
\mathrm{S}\left(\mathrm{~g}_{1}, \mathrm{~g}_{2}\right):=\mathrm{s}_{1} \cdot \mathrm{~g}_{1}-\mathrm{a} \cdot \mathrm{~s}_{2} \cdot \mathrm{~g}_{2}=\mathrm{b} \cdot \mathrm{~s}_{1} \cdot \mathrm{t}_{1}+\mathrm{s}_{1} \cdot \mathrm{R}\left(\mathrm{~g}_{1}\right)-\mathrm{a} \cdot \mathrm{~s}_{2} \cdot \mathrm{R}\left(\mathrm{~g}_{2}\right)
$$

Now B is a Gröbner base iff for every pair of polynomials in $B$ the S-polynomial reduces to 0 w.r.t. $\rightarrow_{B}$.

If $B$ is a Gröbner base for the ideal $I$, then $f \in I$ iff $\left.f{ }^{*}\right]_{B} 0$ and $f \equiv_{I} g$ iff $f$ and $g$ reduce to the same $\rightarrow_{B}$-irreducible element. Thus we can decide ideal membership for $I$, if we have a Gröbner base for I. But a Gröbner base can always be constructed, if a finite set of generators of I ( which always exists by Hilbert's Basis Theorem ) is given.

## (6.4) Buchberger's algorithm.

Let $I$ be an ideal in $\mathbb{Z}\left[X_{1}, \ldots, X_{n}\right]$ and $F$ be a finite set of polynomials such that $\langle F\rangle=I$. As described in (6.3), we can effectively test whether $F$ is a Gröbner base for I. If $F$ is not a Gröbner base, we can extend $F$ by the $\rightarrow_{\mathrm{F}}$-irreducibles of those $S$-polynomials, which do not reduce to 0 , and test again. This completion procedure always terminates with a finite Gröbner base for I ( see e.g. Kandry-Rody-Kapur (1988) for more details ). This termination property is a consequence of Dicksons Lemma ( Dickson (1913)), which holds for free commutative monoids, but not for free monoids ( see e.g. Mora (1985) ).

In the sequel, the following notation will be convenient: Let $h_{1}, \ldots, h_{m}$ be elements of $\mathbb{Z}\left[\mathrm{X}_{1}, \ldots, \mathrm{X}_{\mathrm{n}}\right]$. We denote the $1 \times m$-matrix $\left(\mathrm{h}_{1}, \ldots, \mathrm{~h}_{\mathrm{m}}\right)$ by $\underline{\mathrm{h}}$ and the $\mathrm{m} \times 1$ matrix $\left(\mathrm{h}_{1}, \ldots, \mathrm{~h}_{\mathrm{m}}\right)^{\mathrm{T}}$ ( here ${ }^{\mathrm{T}}$ denotes the transpose of matrices) by $\mid \mathrm{h}$.
Let (*) $\mathrm{f}_{1} \mathrm{x}_{1}+\ldots+\mathrm{f}_{\mathrm{r}} \mathrm{x}_{\mathrm{r}}=\mathrm{f}_{0}$ be an (inhomogeneous) linear equation in $\mathbb{Z}\left[\mathrm{X}_{1}, \ldots, \mathrm{X}_{\mathrm{n}}\right]$. According to Section 3 we have to find one solution for (*) and finitely many generators of the solutions of the homogeneous equation $(* *) \mathrm{f}_{1} \mathrm{x}_{1}+\ldots+\mathrm{f}_{\mathrm{r}} \mathrm{x}_{\mathrm{r}}=0$.
We first construct a Gröbner base $B=\left\{g_{1}, \ldots, g_{s}\right\}$ for $I:=\left\langle\left\{f_{1}, \ldots, f_{r}\right\}\right\rangle$. Since $<B>=I$, there exist an $r \times s$-matrix $P$ and an $s \times r$-matrix $Q$ with entries in $\mathbb{Z}\left[X_{1}, \ldots, X_{n}\right]$ such that $\underline{f} \cdot \mathbf{P}=\mathrm{g}$ and $\mathrm{g} \cdot \mathbf{Q}=\underline{\mathbf{f}}$. This matrices can be obtained as by-products of the Gröbner base construction.
Obviously, (*) has a solution iff $f_{0} \in I$. Hence, if (*) has a solution, then $f_{0}$ reduces to 0 w.r.t $\rightarrow_{B}$. This yields $p_{1}, \ldots, p_{s} \in \mathbb{Z}\left[X_{1}, \ldots, X_{n}\right]$ such that $g \cdot p=f_{0}$. But then $P \cdot l p$ is a solution of (*).
We now assume that we already have finitely many generators $\left|z^{(1)}, \ldots,\right| z^{(L)}$ of the solutions of the equation $(++) g_{1} x_{1}+\ldots+g_{s} x_{s}=0$. Then $P \cdot\left|z^{(1)}, \ldots, P \cdot\right| z^{(L)}$ are solutions of (**), but in general they do not generate all solutions. Let $\mathrm{E}_{\mathrm{r}}$ be the $\mathrm{r} \times \mathrm{r}$ identity matrix and let $\mathrm{It}^{(1)}, \ldots, \mathrm{It}^{(r)}$ be the columns of the matrix $P Q-E_{r}$. Since $\underset{\underline{f}}{\underline{f}}\left(P Q-E_{r}\right)=\underline{f} \cdot P Q-\underline{f} \cdot E_{r}$ $=g \cdot Q-\underline{f}=\underline{0}$, these columns are solutions of $(* *)$.

LEMMA 6.5. The finitely many vectors $\mathrm{P} \cdot\left|\mathrm{z}^{(1)}, \ldots, \mathrm{P} \cdot\right| \mathrm{z}^{(\mathrm{L})},\left|\mathrm{t}^{(1)}, \ldots,\right| \mathrm{t}^{(\mathrm{r})}$ are solutions of $(* *)$ and they generate all solutions of this equation.
PROOF. Let $\mathrm{lq}=\left(\mathrm{q}_{1}, \ldots, \mathrm{q}_{\mathrm{r}}\right)^{\mathrm{T}}$ be an arbitrary solution of $(* *)$. Then $\mathrm{Q} \cdot \mathrm{lq}$ is a solution of $(++)$ and thus there are $a_{1}, \ldots, a_{L} \in \mathbb{Z}\left[X_{1}, \ldots, X_{n}\right]$ such that $Q \cdot\left|q=a_{1} \cdot\right| z^{(1)}+\ldots+a_{L} \mid z^{(L)}$. Now $l q=P Q \cdot\left|q-\left(P Q-E_{r}\right) \cdot\right| q=a_{1} \cdot\left(P \cdot \mid z^{(1)}\right)+\ldots+a_{L} \cdot\left(P \cdot \mid z^{(L)}\right)+q_{1} \cdot \mid t^{(1)}+\ldots+q_{r} \cdot \mathrm{lt}^{(\mathrm{r})}$.

We now show how to solve the equation (++) $g_{1} x_{1}+\ldots+g_{s} x_{s}=0$, if $B=\left\{g_{1}, \ldots, g_{s}\right\}$ is a Gröbner base.
For a set $\left\{\mathrm{q}_{1}, \ldots, \mathrm{q}_{\mathrm{s}}\right\}$ of polynomials the complexity measure $B S\left(q_{1}, \ldots, q_{s}\right)$ is defined as follows: Let $t:=\max \left\{\operatorname{HT}\left(\mathrm{q}_{1}\right), \ldots, \operatorname{HT}\left(\mathrm{q}_{\mathrm{s}}\right)\right.$ and for all $\mathrm{i}, 1 \leq \mathrm{i} \leq \mathrm{s}$, let $\mathrm{a}_{\mathrm{i}}:=\operatorname{coeff}\left(\mathrm{t}, \mathrm{q}_{\mathrm{i}}\right)$ ( Note that $a_{i}=0$ for $\left.H T\left(q_{i}\right)<t\right)$. Then $B S\left(q_{1}, \ldots, q_{s}\right):=\left(\left|a_{1}\right|+\ldots+\left|a_{s}\right|\right) \cdot t$.
Now $t$ is called the term and $\left|a_{1}\right|+\ldots+a_{s} \mid$ the coefficient of $B S\left(q_{1}, \ldots, q_{s}\right)$. We define

$$
a \cdot t=B S\left(q_{1}, \ldots, q_{s}\right)<B S\left(q_{1}^{\prime}, \ldots, q_{s}^{\prime}\right)=a^{\prime} \cdot t^{\prime} \text { iff } t<t^{\prime} \text { or } t=t^{\prime} \text { and } a<a^{\prime}
$$

Let $S\left(g_{i}, g_{j}\right)=s_{i} \cdot g_{i}-a \cdot s_{j} \cdot g_{j}=b \cdot s_{i} \cdot t_{i}+s_{1} \cdot R\left(g_{i}\right)-a \cdot s_{j} \cdot R\left(g_{j}\right)$ be the $S$-polynomial of $g_{i}$ and $g_{j}$ ( see 6.3 ). Since $B$ is a Gröbner base, we have $S\left(g_{i}, g_{j}\right){ }^{*}{ }_{B} 0$. This derivation yields polynomials $w_{1}, \ldots, w_{s}$ such that

$$
S\left(g_{i}, g_{j}\right)=\sum_{k=1}^{k=s} w_{k} \cdot g_{k}
$$

and $\operatorname{BS}\left(w_{1} \cdot g_{1}, \ldots, w_{s} \cdot g_{s}\right)=c \cdot t^{\prime}$ for some $t^{\prime}<s_{i} \cdot t_{i}$, if $b=0$, or $B S\left(w_{1} \cdot g_{1}, \ldots, w_{s} \cdot g_{s}\right)=b \cdot s_{i} \cdot t_{i}$, if $b \neq 0$.
Now $s_{i} \cdot g_{i}-a \cdot s_{j} \cdot g_{j}=S\left(g_{i} \cdot g_{j}\right)=w_{1} \cdot g_{1}+\ldots+w_{s} \cdot g_{s}$ implies $w_{1} \cdot g_{1}+\ldots+\left(w_{i}-s_{i}\right) \cdot g_{i}+\ldots+$ $\left(w_{j}+a \cdot s_{j}\right) \cdot g_{i}+\ldots+w_{s} \cdot g_{s}=0$. Thus $\mid w_{i j}:=\left(w_{1}, \ldots, w_{i}-s_{i}, \ldots, w_{j}+a \cdot s_{j}, \ldots, w_{s}\right)^{T}$ is a solution of the equation ( ++ ).

LEMMA 6.6. The finitely many vectors $1 \mathrm{w}_{\mathrm{ij}}$ generate all solutions of $(++)$.
PROOF. Let $\mathrm{l} \mathrm{p}=\left(\mathrm{p}_{1}, \ldots, \mathrm{p}_{\mathrm{s}}\right)^{\mathrm{T}}$ be a solution of $(++)$ and let $\mathrm{t}=\max \left\{\mathrm{HT}\left(\mathrm{g}_{1} \mathrm{p}_{1}\right), \ldots\right.$, $\mathrm{HT}\left(\mathrm{g}_{s} \mathrm{p}_{\mathrm{s}}\right)$ \}. We prove the lemma by induction on $\mathrm{BS}\left(\mathrm{g}_{1} \mathrm{p}_{1}, \ldots, \mathrm{~g}_{\mathrm{s}} \mathrm{p}_{s}\right)$. Since $\mathrm{g} \cdot \mathrm{p}=0$, there exist $\mathrm{i}, \mathrm{j}$ such that $\mathrm{HT}\left(\mathrm{g}_{\mathrm{i}} \mathrm{p}_{\mathrm{i}}\right)=\mathrm{t}=\mathrm{HT}\left(\mathrm{g}_{\mathrm{j}} \mathrm{p}_{\mathrm{j}}\right)$ and $\mathrm{HC}\left(\mathrm{g}_{\mathrm{i}} \mathrm{p}_{\mathrm{i}}\right)$ and $\mathrm{HC}\left(\mathrm{g}_{\mathrm{j}} \mathrm{p}_{\mathrm{j}}\right)$ have different sign. Without loss of generality we assume that $c_{i}:=\mathrm{HC}\left(\mathrm{g}_{\mathrm{i}}\right) \geq \mathrm{HC}\left(\mathrm{g}_{\mathrm{j}}\right)=: \mathrm{c}_{\mathrm{j}}>0$. Obviously, $\mathrm{t}_{\mathrm{i}}$ $:=H T\left(g_{i}\right)$ and $t_{j}:=H T\left(g_{j}\right)$ are divisors of $t$ and thus $\operatorname{lcm}\left(t_{i} t_{j}\right)=s_{i} t_{i}=s_{j} t_{j}$ divides $t$, i.e. there exists $r$ with $r s_{i} t_{i}=r s_{j} t_{j}=t$.
We now consider the case $\mathrm{HC}\left(\mathrm{g}_{\mathrm{i}} \mathrm{p}_{\mathrm{i}}\right)>0$ and $\mathrm{HC}\left(\mathrm{g}_{\mathrm{j}} \mathrm{p}_{\mathrm{j}}\right)<0$ (the other case is similar ). The vector $\mathrm{lq}=\left(\mathrm{q}_{1}, \ldots, \mathrm{q}_{s}\right)^{\mathrm{T}}:=\mathrm{lp}+\mathrm{r} \cdot \mathrm{w}_{\mathrm{ij}}$ is a solution of $(++)$ and we have $\mathrm{g}_{1} \mathrm{q}_{1}=\mathrm{g}_{1} \mathrm{p}_{1}+$ $g_{1} r w_{1}, \ldots, g_{i} q_{i}=g_{i} p_{i}+g_{i} r w_{i}-g_{i} r s_{i}, \ldots, g_{j} q_{j}=g_{j} p_{j}+g_{j} r w_{j}+g_{j} a r s_{j}, \ldots, g_{s} q_{s}=g_{s} p_{s}+g_{s} r w_{s}$. If $\max \left\{\mathrm{HT}\left(\mathrm{g}_{1} \mathrm{q}_{1}\right), \ldots, \mathrm{HT}\left(\mathrm{g}_{\mathrm{s}} \mathrm{q}_{s}\right)\right\}<\mathrm{t}$, the lemma is proved by induction, since the term of BS has decreased. Otherwise $\max \left\{H T\left(g_{1} q_{1}\right), \ldots, H T\left(g_{s} q_{s}\right)\right\}=t$ and we have to calcu-
late the coefficient of $\operatorname{BS}\left(g_{1} q_{1}, \ldots, g_{s} q_{s}\right)$. The triangle inequality yields

$$
B S\left(g_{1} q_{1}, \ldots, g_{s} q_{s}\right) \leq B S\left(g_{1} p_{1}, \ldots, g_{i} p_{i}-g_{i} r_{i}, \ldots, g_{j} p_{j}+g_{j} a s_{j}, \ldots, g_{s} p_{s}\right)+b \cdot t,
$$

since $\mathrm{BS}\left(\mathrm{g}_{1} \mathrm{rw}_{1} \ldots, \mathrm{~g}_{\mathrm{s}} \mathrm{rw} \mathrm{w}_{\mathrm{s}}\right)=\mathrm{r} \cdot \mathrm{b} \cdot \mathrm{s}_{\mathrm{i}} \cdot \mathrm{t}_{\mathrm{i}}=\mathrm{b} \cdot \mathrm{t}\left(\right.$ for $\mathrm{b}>0$ ) or $\mathrm{BS}\left(\mathrm{g}_{1} \mathrm{rw}_{1}, \ldots, \mathrm{~g}_{\mathrm{s}} \mathrm{rw} \mathrm{w}_{\mathrm{s}}\right)$ has a term which is smaller than $t$ for $b=0$ ).
We have $\operatorname{lcoeff}\left(t, g_{i} p_{i}-g_{i} r s_{i}\right)\left|=\left|\operatorname{coeff}\left(t, g_{i} p_{i}\right)\right|-c_{i}\left(\right.\right.$ since $\left.\operatorname{coeff}\left(t, g_{i} p_{i}\right)=H C\left(g_{i} p_{i}\right) \geq c_{i} \geq 0\right)$ and $\left|\operatorname{coeff}\left(\mathrm{t}, \mathrm{g}_{\mathrm{j}} \mathrm{p}_{\mathrm{j}}+\mathrm{g}_{\mathrm{j}} \operatorname{ars} \mathrm{s}_{\mathrm{j}}\right)\right|<\left|\operatorname{coeff}\left(\mathrm{t}, \mathrm{g}_{\mathrm{j}} \mathrm{p}_{\mathrm{j}}\right)\right|+\mathrm{ac}_{\mathrm{j}}\left(\right.$ since $\operatorname{coeff}\left(\mathrm{t}, \mathrm{g}_{\mathrm{j}} \mathrm{p}_{\mathrm{j}}\right)=\mathrm{HC}\left(\mathrm{g}_{\mathrm{j}} \mathrm{p}_{\mathrm{j}}\right)<0$ and $\operatorname{coeff}\left(\mathrm{t}, \mathrm{g}_{\mathrm{j}} \mathrm{ars}_{\mathrm{j}}\right)=\mathrm{ac}_{\mathrm{j}}>0$ ).
Thus BS $\left(g_{1} p_{1}, \ldots, g_{i} p_{i}-g_{i} r_{i}, \ldots, g_{j} p_{j}+g_{j} \operatorname{ars}_{j} \ldots, g_{s} p_{s}\right)<B S\left(g_{1} p_{1}, \ldots, g_{s} p_{s}\right)+\left(a c_{j}-c_{i}\right) \cdot t$ and, since $c_{i}=a \cdot c_{j}+b, B S\left(g_{1} q_{1}, \ldots, g_{s} q_{s}\right)<B S\left(g_{1} p_{1}, \ldots, g_{s} p_{s}\right)$. This completes the proof of Lemma 6.6 by induction on $B S$.

Now we have completely described a method to solve linear equations in $\mathbb{Z}\left[X_{1}, \ldots, X_{n}\right]$.

EXAMPLE 6.7. As an example, consider the equation $f_{1} x_{1}+f_{2} x_{2}+f_{3} x_{3}=f_{0}$ for $f_{0}=$ $X^{3} Y Z^{2}-X^{3} Y^{3} Z^{2}, f_{1}=X^{3} Y Z-X Z^{2}, f_{2}=X Y^{2} Z-X Y Z$ and $f_{3}=X^{2} Y^{2}-Z$.
First, we have to calculate a Gröbner base for the Ideal I, generated by $f_{1}, f_{2}$ and $f_{3}$. Let < be the admissible ordering defined by the matrix

$$
M_{<}=\left(\begin{array}{lll}
1 & 0 & 0 \\
1 & 0 & 1 \\
1 & 1 & 0
\end{array}\right) \text { (that means: first order by total degree and, within a given degree, }
$$

With respect to this ordering, the Buchberger algorithm yields the Gröbner base $B=\left\{g_{1}\right.$, $\left.\mathrm{g}_{2}, \mathrm{~g}_{3}, \mathrm{~g}_{4}, \mathrm{~g}_{5}\right\}$, where $\mathrm{g}_{1}=\mathrm{f}_{2}, \mathrm{~g}_{2}=\mathrm{f}_{3}, \mathrm{~g}_{3}=\mathrm{X}^{2} \mathrm{YZ}-\mathrm{Z}^{2}, \mathrm{~g}_{4}=\mathrm{YZ}^{2}-\mathrm{Z}^{2}$ and $\mathrm{g}_{5}=\mathrm{X}^{2} Z^{2}-Z^{3}$. By keeping track of how the $g_{i}$ are generated in this process, we obtain the transformation matrix $P$ such that $\underline{f} \cdot P=g$ and, by reduction of the $f_{j}$ w.r.t. $\rightarrow_{B}$, we get the matrix $Q$ such that $\mathrm{g} \cdot \mathrm{Q}=\underline{\mathrm{f}}$. In our example
$P=\left(\begin{array}{lllll}0 & 0 & 0 & 0 & 0 \\ 1 & 0 & -X & X Y & -Z X-X^{3} Y \\ 0 & 1 & Z & -Y Z+Z & Z^{2}+X^{2} Y Z-X^{2} Z\end{array}\right)$ and $Q=\left(\begin{array}{ccc}0 & 1 & 0 \\ 0 & 0 & 1 \\ X & 0 & 0 \\ 0 & 0 & 0 \\ 0 & 0 & 0\end{array}\right)$.
We now determine whether $f_{0} \in I=\langle B\rangle$, i.e. whether $f_{0}$ reduces to 0 w.r.t. $\rightarrow_{B}$ : $\mathrm{f}_{0} \rightarrow_{\mathrm{B}} \mathrm{f}_{0}-\mathrm{g}_{5} \cdot X Y=X Y Z^{3}-X^{3} Y^{3} Z^{2} \rightarrow_{\mathrm{B}} \mathrm{f}_{0}-\mathrm{g}_{5} \cdot X Y+\mathrm{g}_{3} \cdot X Y^{2} Z=X Y Z^{3}-X Y^{2} Z^{3} \rightarrow_{B}$ $\mathrm{f}_{0}-\mathrm{g}_{5} \cdot X Y+\mathrm{g}_{3} \cdot X Y^{2} Z+\mathrm{g}_{4} \cdot X Y Z=X Y Z^{3}-X Y Z^{3}=0$.
Thus $f_{0}=g_{1} \cdot 0+g_{2} \cdot 0+g_{3} \cdot\left(-X Y^{2} Z\right)+g_{4} \cdot(-X Y Z)+g_{5} \cdot X Y \in\langle B\rangle=I$ and we can use the transformation matrix $P$ to obtain a solution of the equation $f_{1} x_{1}+f_{2} x_{2}+f_{3} x_{3}=f_{0}$ :

$$
P \cdot\left(0,0,-X Y^{2} Z,-X Y Z, X Y\right)^{T}=\left(0,-X^{2} Y Z-X^{4} Y^{2}, X^{3} Y^{2} Z-X^{3} Y Z\right)^{T}
$$

The next step is to determine the solutions $l w_{i j}$ of the equation $g_{1} x_{1}+\ldots+g_{5} x_{5}=0$.
$S\left(g_{1}, g_{2}\right)=g_{1} \cdot X-g_{2} \cdot Z=-X^{2} Y Z+Z^{2}=-g_{3}$ and thus $g_{1} \cdot(-X)+g_{2} \cdot Z+g_{3} \cdot(-1)+g_{4} \cdot 0+$ $g_{5} \cdot 0=0$. That means ${ }^{\prime} w_{1,2}=(-X, Z,-1,0,0)^{T}$.
$S\left(g_{1}, g_{3}\right)=g_{1} \cdot X-g_{3} \cdot Y=-X^{2} Y Z+Y Z^{2}=-g_{3}-Z^{2}+Y Z^{2}=-g_{3}+g_{4}$ and thus we get $\mathrm{lw}_{1,3}=(-X, 0, Y-1,1,0)^{T}$.
Similar computations yield the other vectors $\mid \mathrm{w}_{\mathrm{ij}}$ :
$\mathrm{Iw}_{1,4}=(-Z, 0,0, X Y, 0)^{T}, \quad \quad \mathrm{lw}_{1,5}=\left(-X Y, 0,-Z, Y Z+Z, Y^{2}\right)^{T}$,
$l w_{2,3}=(0,-Z, Y, 1,0)^{T}, \quad \quad l w_{2,4}=\left(0,-Z^{2}, Z, X^{2} Y, 0\right)^{T}$,
${ }_{l w_{2,5}}=\left(0,-Z^{2}, 0, Y Z+Z, Y^{2}\right)^{T}, \quad l w_{3,4}=\left(0,0,-Z, X^{2}, 1\right)^{T}$,
$\mathrm{lw}_{3,5}=(0,0,-\mathrm{Z}, \mathrm{Z}, \mathrm{Y})^{\mathrm{T}}, \quad \quad \quad \mathrm{w}_{4,5}=\left(0,0,0,-\mathrm{X}^{2}+\mathrm{Z}, \mathrm{Y}-1\right)^{\mathrm{T}}$.
Now we use the transformation matrix $P$ to obtain solutions of the homogeneous equation $\mathrm{f}_{1} \mathrm{x}_{1}+\mathrm{f}_{2} \mathrm{x}_{2}+\mathrm{f}_{3} \mathrm{x}_{3}=0$ :
$\mathrm{P} \cdot \mathrm{l} \mathrm{w}_{1,2}=(0,0,0)^{\mathrm{T}}$,

$$
P \cdot \mid w_{1,4}=\left(0, X^{2} Y^{2}-Z,-X Y^{2} Z+X Y Z\right)^{T}
$$

$$
P \cdot \mid w_{2,3}=(0,0,0)^{T}
$$

$$
P \cdot\left|w_{2,5}=P \cdot\right| w_{1,5}=(-X Y) \cdot P \cdot \mid w_{1,4}
$$

$$
P \cdot\left|w_{3,5}=-P \cdot\right| w_{2,4}=(-X) \cdot P \cdot \mid w_{1,4}
$$

$$
\begin{aligned}
& \mathrm{P} \cdot \mid \mathrm{w}_{1,3}=(0,0,0)^{\mathrm{T}} \\
& \mathrm{P} \cdot\left|\mathrm{w}_{1,5}=(-X Y) \cdot \mathrm{P} \cdot\right| \mathrm{w}_{1,4} \\
& \mathrm{P} \cdot\left|\mathrm{w}_{2,4}=\mathrm{X} \cdot \mathrm{P} \cdot\right| \mathrm{w}_{1,4}, \\
& \mathrm{P} \cdot \mid \mathrm{w}_{3,4}=(0,0,0)^{\mathrm{T}}, \\
& \mathrm{P} \cdot\left|\mathrm{w}_{4,5}=\mathrm{P} \cdot\right| \mathrm{w}_{3,5}=(-X) \cdot \mathrm{P} \cdot \mid \mathrm{w}_{1,4}
\end{aligned}
$$

The solution $P \cdot \mid w_{1,4}=\left(0, X^{2} Y^{2}-Z,-X Y^{2} Z+X Y Z\right)^{T}$ thus obtained does not generate all solutions of $f_{1} x_{1}+f_{2} x_{2}+f_{3} x_{3}=0$. In addition, we need the columns of the matrix

$$
P \cdot Q-E_{3}=\left(\begin{array}{lll}
-1 & 0 & 0 \\
-X^{2} & 0 & 0 \\
X Z & 0 & 0
\end{array}\right)
$$

All solutions of the homogeneous equation $f_{1} x_{1}+f_{2} x_{2}+f_{3} x_{3}=0$ are generated by the two solutions $\left(0, X^{2} Y^{2}-Z,-X Y^{2} Z+X Y Z\right)^{T}$ and $\left(-1,-X^{2}, X Z\right)^{T}$.

EXAMPLE 6.8. As a second example, we consider the equation $\mathrm{Xx}_{1}+\mathrm{Xx}_{2}=\mathrm{x}_{2}+\mathrm{X}^{2} \mathrm{x}_{3}$ of Section 4, but now we want to solve it in $\mathbb{Z}[X]$. Hence we have to solve the homogeneous equation $f_{1} x_{1}+f_{2} x_{2}+f_{3} x_{3}=0$ for $f_{1}=X, f_{2}=X-1$ and $f_{3}=-X^{2}$. It is easy to see that $\left\langle\left\{f_{1}, f_{2}, f_{3}\right\}\right\rangle=\mathbb{Z}[X]$ and that $B=\left\{g_{1}\right\}$ for $g_{1}=1$ is the corresponding Gröbner base. The transformation matrices are $P=(1,-1,0)^{T}$ and $Q=\left(X, X-1,-X^{2}\right)$.
Obviously, the equation $g_{1} x_{1}$ has only the trivial solution $x_{1}=0$. Thus the columns of

$$
P \cdot Q-E_{3}=\left(\begin{array}{ccc}
X-1 & X-1 & -X^{2} \\
-X & -X & X^{2} \\
0 & 0 & -1
\end{array}\right)
$$

i.e. $(X-1,-X, 0)^{T}$ and $\left(-X^{2}, X^{2},-1\right)^{T}$, generate all solutions of $X x_{1}+X_{2}=x_{2}+X^{2} x_{3}$ in $(\mathbb{Z}[X])^{3}$.

## 7. AGnH-Unification

It is easy to see that $\mathrm{S}(\mathrm{AGnH})$ is isomorphic to the ring $\mathbb{Z}<\mathrm{X}_{1}, \ldots, X_{\mathrm{n}}>$, i.e. the polynomial ring over $\mathbb{Z}$ in the non-commuting indeterminates $X_{1}, \ldots, X_{n}$. Unfortunately, for $n \geq 2$ this ring is not noetherian ( see Mora (1985) ) and the membership problem for finitely generated two-sided ideals is undecidable ( Kandry-Rody-Weispfenning (1988) ). Fortunately, we are not interested in two-sided ideals, but only in right ideal. The solutions of a homogeneous equation $f_{1} x_{1}+\ldots+f_{r} x_{r}=0$ are only closed under right multiplication and the inhomogeneous equation $f_{1} x_{1}+\ldots+f_{T} x_{T}=f_{0}$ has a solution iff $f_{0}$ is a member of the right ideal generated by $f_{1}, \ldots ., f_{r}$. Though, for $n \geq 2, \mathbb{Z}<X_{1}, \ldots, X_{n}>$ is not even right noetherian (i.e. there are right ideals in $\mathbb{Z}<X_{1}, \ldots, X_{n}>$, which are not finitely generated), the set of solutions of a homogeneous equation $f_{1} x_{1}+\ldots+f_{r} x_{T}=0$ is a finitely generated right $\mathbb{Z}<X_{1}, \ldots, X_{\mathrm{n}}>$-semimodule and the membership problem for finitely generated right ideals is decidable in $\mathbb{Z}<X_{1}, \ldots, X_{n}>$ ( see Section 8 and 9 ). This yields

PROPOSITION 7.1. For any $\mathrm{n} \geq 0$ the theory AGnH is unitary and it is also unitary w.r.t. unification with constants.

## 8. "Gröbner bases" for finitely generated right ideals in $\mathbb{Z}<X_{1}, \ldots, X_{n}>$

The construction of Gröbner bases for finitely generated right ideals in $K<X_{1}, \ldots, X_{n}>$, where K is a field, is very easy ( Mora (1985)). For $\mathbb{Z}<X_{1}, \ldots, X_{n}>$ one has to be more careful.
The rôle of terms in the commutative case is now played by words over the alphabet $\Sigma_{\mathrm{n}}$ $:=\left\{X_{1}, \ldots, X_{n}\right\}$. Let $W_{n}$ be the set of these words, i.e. the free monoid generated by $\Sigma_{n}$, and let 1 denote the empty word.
A total ordering < on $\mathrm{W}_{\mathrm{n}}$ is called admissible iff the following two conditions hold:
(1) For all $\mathrm{s}, \mathrm{t}, \mathrm{r} \in \mathrm{W}_{\mathrm{n}} \mathrm{s}<\mathrm{t}$ implies $\mathrm{sr}<\mathrm{tr}$ ( compatibility with right concatenation).
(2) For all $s \in W_{n}$ the set $\left\{t \in W_{n} ; t<s\right\}$ is finite.

LEMMA 8.1. Let < be an admissible ordering on $\mathrm{W}_{\mathrm{n}}$.
(1) < is order-isomorphic to $\omega$ and thus noetherian.
(2) $1<t$ for all $t \in W_{n} \backslash\{1\}$.
(3) $\mathrm{s}=$ tr for $\mathrm{r} \neq 1$ implies $\mathrm{s}>\mathrm{t}$.

Examples of admissible orderings are graded lexicographical orderings and, more general,
all suffle-compatible total orders ( see Leeb-Pirillo (1988)). The complete characterization of all concatenation-compatible ( resp. right concatenation-compatible) linear orderings is still an open problem.
We now extend admissible orderings to monomials and polynomials.
DEFINITION 8.2. Let $<$ be an admissible ordering on $W_{n}$.
(1) Let $a, b \in \mathbb{Z}$ and $s, t \in W_{n}$. Then as $<b t$ iff $s<t$ or $s=t$ and $l a l<|b|$ or $s=t$ and $|a|=|\mathrm{b}|$ and $\mathrm{a}<\mathrm{b}$. This defines a well-ordering on the monomials of $\left.\mathbb{Z}<\mathrm{X}_{1}, \ldots, \mathrm{X}_{\mathrm{n}}\right\rangle$.
(2) Let $f=\Sigma a_{i} s_{i}$ and $g=\Sigma b_{i} t_{i}$ be two polynomials, i.e. elements of $\mathbb{Z}\left\langle X_{1}, \ldots, X_{n}\right\rangle$. Then we define $f<g$ iff $\left\{\ldots a_{i} s_{i}, \ldots\right\} \ll\left\{\ldots b_{i} t_{i}, \ldots\right\}$, where $\ll$ denotes the multiset ordering ( see Dershowitz-Manna (1979) ) induced by the ordering <on monomials.
(3) Let $f$ be a polynomial. We write $f=$ at $+R(f)$ if $t$ is the maximal (w.r.t. <) word in $f$ ( $t=H W(f)$ ) and $a$ is the coefficient of $\operatorname{tin} f(a=H C(f)$ ).
(4) For a set $F$ of polynomials in $\mathbb{Z}<X_{1}, \ldots, X_{n}>$, the reduction relation $\rightarrow_{F}$ is defined as in Section 6, 6.3.

For $\mathrm{K}<\mathrm{X}_{1}, \ldots, \mathrm{X}_{\mathrm{n}}>$, Mora (1985) has described a very easy algorithm, which transforms a finite set F of polynomials into a "Gröbner base" ( see Mora (1985) for the definition of Gröbner bases in this case ):
Start with $\mathrm{F}_{0}:=\mathrm{F}$. As long as there are polynomials $f, g$ in $\mathrm{F}_{\mathrm{k}}$, such that $H W(f)$ is a prefix of $H W(g)$, $g$ can be reduced by $f$ to a smaller polynomial $g$. Define $F_{k+1}=\left(F_{k} \backslash\{g\}\right)$ $\cup\left\{g^{\prime}\right\}$ and continue with $F_{k+1}$ in place of $F_{k}$
This process terminates after finitely many steps and yields a finite set $G$ of polynomials, which generates the same right ideal as $F$ and has the following property:
For two different elements $f$ and $g$ of $G, H W(f)$ and $H W(g)$ are not comparable w.r.t. the prefix-ordering ( i.e. for any word $r, H W(f) \cdot r \neq H W(g)$ and $H W(g) \cdot r \neq H W(f)$ ).
For $\mathbb{Z}<X_{1}, \ldots, X_{n}>$, we encounter the following problem: For $f=a \cdot t+R(f)$ and $g=b \cdot t \cdot r+$ $R(g)$ with $t, r \in W_{n}, a, b \in \mathbb{Z}$ and $|a|>|b|, H W(f)$ is prefix of $H W(g)$, but the head monomial of $g$ can not be reduced by $f$. If, in addition, $b$ devides $a$, it may become necessary to increase the actual set of polynomials ( see Case 4 below ). Since Dickson's Lemma does not hold for free monoids, we have to be very careful, if we want to obtain a terminating algorithm.

ALGORITHM 8.3. This is the informal describtion of an algorithm, which transforms a finite set of polynomials $\left\{f_{1}, \ldots, f_{m}\right\}$ into a "Gröbner base", which defines the same right ideal.
In the beginning, $\mathrm{F}_{0}:=\left\{\mathrm{f}_{1}, \ldots, \mathrm{f}_{\mathrm{m}}\right\}$ and all pairs of indices are unmarked.
Assume that $F_{k}(k \geq 0)$ is already defined. If there is the zero polynomial 0 in $F_{k}$, we erase it. While there are $f:=f_{i}$ and $g:=f_{j}$ in $F_{k}$ such that
(1) $(i, j)$ is not marked and
(2) $f=a \cdot t+R(f)$ and $g=b \cdot \operatorname{tr}+R(g)$ for some $a, b \in \mathbb{Z}$ and $t, r \in W_{n}$,
we do the following:
Case 1: $\mathrm{r}=1$.
Without loss of generality we may assume that lal $\geq \mathrm{lbl}$. Let $\mathrm{a}=\mathrm{bc}+\mathrm{d}$ for some $\mathrm{c}, \mathrm{d}$ such that $0 \leq \mathrm{d}<\mathrm{lbl} \leq$ lal.
Define $f_{1}:=f-g \cdot c=d \cdot t+R(f)-R(g) \cdot c$ and $F_{k+1}:=\left(F_{k} \backslash\{f\}\right) \cup\left\{f_{1}\right\}$. We do not have to mark ( $i, j$ ), since $f=f_{i}$ is removed.
Obviously, $\mathrm{f}_{1}<\mathrm{f}$ and $\mathrm{f}=\mathrm{f}_{1}+\mathrm{g} \cdot \mathrm{c}$. Hence $\mathrm{F}_{\mathrm{k}+1}$ generates the same right ideal as $\mathrm{F}_{\mathrm{k}}$, but f is replaced by the smaller polynomial $f_{1}$.

Case 2. $\mathrm{r} \neq 1$ and $\mathrm{lal} \leq \mathrm{l} \mathrm{b}$.
Let $\mathrm{b}=\mathrm{ac}+\mathrm{d}$ for some $\mathrm{c}, \mathrm{d}$ such that $0 \leq \mathrm{d}<|\mathrm{al} \leq| \mathrm{bl}$.
Define $\mathrm{g}_{1}:=\mathrm{g}-\mathrm{f} \cdot \mathrm{cr}=\mathrm{d} \cdot \mathrm{tr}+\mathrm{R}(\mathrm{g})-\mathrm{R}(\mathrm{f}) \cdot \mathrm{cr}$ and $\mathrm{F}_{\mathrm{k}+1}:=\left(\mathrm{F}_{\mathrm{k}} \backslash\{\mathrm{g}\}\right) \cup\left\{\mathrm{g}_{1}\right\}$.
Obviously, $g_{1}<g$, and $g=g_{1}+f \cdot c r$. Hence $F_{k+1}$ generates the same right ideal as $F_{k}$, but g is replaced by the smaller polynomial $\mathrm{g}_{1}$.
Case 3. $\mathrm{r} \neq 1$, $\mathrm{lal}>\mathrm{lb\mid}$ and $\mathrm{lb} \mid$ does not devide lal.
Let $\mathrm{a}=\mathrm{bc}+\mathrm{d}$ for some $\mathrm{c}, \mathrm{d}$ such that $0<\mathrm{d}<\mathrm{lb} \mid<\mathrm{lal}$. We define $\mathrm{g}_{1}:=\mathrm{f} \cdot \mathrm{r}-\mathrm{g} \cdot \mathrm{c}=\mathrm{d} \cdot \mathrm{tr}+$ $R(f) \cdot r-R(g) \cdot c$. Since the words occurring in $R(f) \cdot r$ and $R(g) \cdot c$ are smaller than $t r$, we have $\mathrm{HW}\left(\mathrm{g}_{1}\right)=\operatorname{tr}, \mathrm{HC}\left(\mathrm{g}_{1}\right)=\mathrm{d}$ and $\mathrm{R}\left(\mathrm{g}_{1}\right)=\mathrm{R}(\mathrm{f}) \cdot \mathrm{r}-\mathrm{R}(\mathrm{g}) \cdot \mathrm{c}$. Obviously, $\mathrm{g}_{1}<\mathrm{g}, \mathrm{g}_{1} \in<\mathrm{F}_{\mathrm{k}}>$ and the pair $g_{1}$, $g$ satisfies Case 1 . Hence we define $g_{2}:=g-g_{1} \cdot c_{1}$ (where $b=d c_{1}+d_{1}$, $\left.0 \leq d_{1}<d\right)$ and $F_{k+1}:=\left(F_{k} \backslash\{g\}\right) \cup\left\{g_{1}, g_{2}\right\}$. Since $g_{1}, g_{2}<g$ and $g=g_{2}+g_{1} \cdot c, F_{k+1}$ generates the same right ideal as $\mathrm{F}_{\mathrm{k}}$, but g is replaced by the two smaller polynomials $\mathrm{g}_{1}$ and $\mathrm{g}_{2}$.
Case 4. $\mathrm{r} \neq 1,|\mathrm{la\mid}>| \mathrm{bl}$ and $\mid \mathrm{bl}$ devides lal, i.e. there exists c such that $\mathrm{a}=\mathrm{bc}$.
Define $g_{1}:=\mathrm{f} \cdot \mathrm{r}-\mathrm{g} \cdot \mathrm{c}=\mathrm{R}(\mathrm{f}) \cdot \mathrm{r}-\mathrm{R}(\mathrm{g}) \cdot \mathrm{c}$. Now $\mathrm{g}_{1}<\mathrm{g}$, but since $\mathrm{Icl} \neq 1$, g can not be represented using $\mathrm{g}_{1}$. We distinguish the following cases:
Case 4.1. There is $h \in \cup_{i \leq k} F_{i}$ with the property $\mathrm{HW}\left(\mathrm{g}_{1}\right)=\mathrm{HW}(\mathrm{h})$.
Case 4.1.1. $\mathrm{h} \in \mathrm{F}_{\mathrm{k}}$ and $\left|\mathrm{HC}\left(\mathrm{g}_{1}\right)\right|<|\mathrm{HC}(\mathrm{h})|$.
We have $g_{1}<h$ and $h$ may be reduced by $g_{1}$ to some $h_{1}<h($ see Case 1 ).
Define $\mathrm{F}_{\mathrm{k}+1}:=\left(\mathrm{F}_{\mathrm{k}} \backslash\{\mathrm{h}\}\right) \cup\left\{\mathrm{g}_{1}, \mathrm{~h}_{1}\right\}$ and mark $(\mathrm{i}, \mathrm{j}) . \mathrm{F}_{\mathrm{k}+1}$ generates the same right ideal as $F_{k}$, but $h$ is replaced by the two smaller polynomials $g_{1}$ and $h_{1}$.
Case 4.1.2. $\mathrm{h} \in \mathrm{F}_{\mathrm{k}}$ and $\left|\mathrm{HC}\left(\mathrm{g}_{1}\right)\right| \geq|\mathrm{HC}(\mathrm{h})|$.
Then $g_{1}$ may be reduced by $h$ to a smaller polynomial $g_{2}$ ( see Case 1 ). If $g_{2}=0, F_{k+1}:=$ $F_{k}$ and we mark $(i, j)$. Otherwise we continue with $g_{2}$ in place of $g_{1}$.
Case 4.1.3. $\mathrm{h} \notin \mathrm{F}_{\mathrm{k}}$ and there is no polynomial in $\mathrm{F}_{\mathrm{k}}$ which has $\mathrm{HW}(\mathrm{h})$ as head word.
Thus the monomial $\mathrm{HC}(\mathrm{h}) \mathrm{HW}(\mathrm{h})$ has been reduced in some previous step. It is easy to see, that then $\mathrm{HC}(\mathrm{h}) \mathrm{HW}(\mathrm{h})$ can also be reduced by $\rightarrow_{\mathrm{F}_{\mathrm{k}}}$. If we have $\left|\mathrm{HC}\left(\mathrm{g}_{1}\right)\right| \geq|\mathrm{HC}(\mathrm{h})|$, $\mathrm{g}_{1}$ can be reduced and we proceed as in Case 4.1.2.

Otherwise, i.e. if $\left|\mathrm{HC}\left(\mathrm{g}_{1}\right)\right|<|\mathrm{HC}(\mathrm{h})|$, we define $\mathrm{F}_{\mathrm{k}+1}:=\mathrm{F}_{\mathrm{k}} \cup\left\{\mathrm{g}_{1}\right\}$ and mark $(\mathrm{i}, \mathrm{j})$.
Case 4.2. There is no $h \in \cup_{i \leq k} F_{i}$ with the property $H W\left(g_{1}\right)=H W(h)$.
In this case we also define $F_{k+1}:=F_{k} \cup\left\{g_{1}\right\}$ and mark $(i, j)$.

This completes the description of Algorithm 8.3. We shall soon show that this algorithm always terminates with a finite set of polynomials $G$, whose properties justify the name Gröbner base. But first, we consider an example.

EXAMPLE 8.4. Let $f_{1}=2 a b c-b c, f_{2}=3 a b-2 b, f_{3}=5 a b d-b c$ and $f_{4}=b c-5 b d$ be polynomials in $\mathbb{Z}<a, b, c, d>$. We take the graded lexicograpical ordering with $a>b>c>d$ as admissible ordering (i.e. $\mathbf{u}<\mathbf{v}$ iff $|\mathbf{u}|<\mid \mathbf{v |}$ or $\left|\mathbf{u l}=|\mathbf{v}|\right.$ and $\mathbf{u}<_{\mathrm{lex}} \mathrm{v}$ ) and run Algorithm 8.3 with input $F_{0}:=\left\{f_{1}, f_{2}, f_{3}, f_{4}\right\}$.

1) For $f_{1}$ and $f_{2}$ we have Case 3.

Define $f_{5}:=f_{2} \cdot c-f_{1}=a b c-b c$ and $f_{6}:=f_{1}-f_{5} \cdot 2=b c$. Now $f_{1}$ is replaced by $f_{5}, f_{6}$, which yields $F_{1}=\left\{f_{2}, f_{3}, f_{4}, f_{5}, f_{6}\right\}$. We have $f_{1}=f_{5} \cdot 2+f_{6}$.
2) For $f_{2}$ and $f_{3}$ we have Case 2 .

Define $f_{7}:=f_{3}-f_{2} \cdot d=2 a b d-b c+2 b d$ and replace $f_{3}$ by $f_{7}$, which yields $F_{2}=\left\{f_{2}, f_{4}, f_{5}, f_{6}\right.$, $\left.f_{7}\right\}$. We have $f_{3}=f_{7}+f_{2} \cdot d$.
3) For $f_{2}$ and $f_{5}$ we have Case 4.

Define $f_{8}=f_{2} \cdot c-f_{5} \cdot 3=b c=f_{6}$. Hence we have Case 4.1.2 and since $f_{6}$ reduces $f_{8}$ to $0, F_{3}$ $=F_{2}=\left\{f_{2}, f_{4}, f_{5}, f_{6}, f_{7}\right\}$ and the index pair (2,5) is marked.
4) For $f_{2}$ and $f_{7}$ we have Case 3.

Define $f_{9}:=f_{2} \cdot d-f_{7}=a b d-4 b d+b c$ and $f_{10}=f_{7}-f_{9} \cdot 2=-3 b c+10 b d$. Now $f_{7}$ is replaced by $\mathrm{f}_{9}$ and $\mathrm{f}_{10}$, which yields $\mathrm{F}_{4}=\left\{\mathrm{f}_{2}, \mathrm{f}_{4}, \mathrm{f}_{5}, \mathrm{f}_{6}, \mathrm{f}_{9}, \mathrm{f}_{10}\right\}$. We have $\mathrm{f}_{7}=\mathrm{f}_{10}+\mathrm{f}_{9} \cdot 2$.
5) For $f_{2}$ and $f_{9}$ we have Case 4.

Define $f_{11}:=f_{2} \cdot d-f_{9} \cdot 3=-3 b c+10 b d$. Now $H W\left(f_{11}\right)=H W\left(f_{4}\right)$ and $f_{4}$ reduces $f_{11}$ to the polynomial $f_{12}:=f_{11}+f_{4} \cdot 3=-5 b d$ ( Case 4.1.2 ). We continue with $f_{12}$ in place of $f_{11}$ and have Case 4.2, since bd has not yet occured as head word. Hence $F_{5}:=F_{4} \cup\left\{f_{12}\right\}$ and $(2,5)$ and $(2,9)$ are already marked.
6) For $f_{4}$ and $f_{6}$ we have Case 1.

Define $\mathrm{f}_{13}:=\mathrm{f}_{4}-\mathrm{f}_{6}=\mathrm{f}_{12}$ and $\mathrm{F}_{6}:=\mathrm{F}_{5} \backslash\left\{\mathrm{f}_{4}\right\}=\left\{\mathrm{f}_{2}, \mathrm{f}_{5}, \mathrm{f}_{6}, \mathrm{f}_{9}, \mathrm{f}_{10}, \mathrm{f}_{12}\right\}$.
7) For $f_{6}$ and $f_{10}$ we have Case 1.

Define $\mathrm{f}_{14}:=\mathrm{f}_{10}+\mathrm{f}_{6} \cdot 3=10 \mathrm{bd}$ and $\mathrm{F}_{7}:=\left\{\mathrm{f}_{2}, \mathrm{f}_{5}, \mathrm{f}_{6}, \mathrm{f}_{9}, \mathrm{f}_{12}, \mathrm{f}_{14}\right\}$.
8) For $f_{12}$ and $f_{14}$ we have Case 1 .

Since $f_{14}=f_{12} \cdot(-2), f_{14}$ can be eliminated and we get $F_{8}=\left\{f_{2}, f_{5}, f_{6}, f_{9}, f_{12}\right\}$, where $(2,5)$ and $(2,9)$ are marked.

Hence Algorithm 8.3 terminates with $G:=F_{8}=\left\{f_{2}, f_{5}, f_{6}, f_{9}, f_{12}\right\}$. The elements of $G$ are $g_{1}:=f_{2}=3 a b-2 b, g_{2}:=f_{5}=a b c-b c, g_{3}:=f_{6}=b c, g_{4}:=f_{9}=a b d-4 b d+b c$ and $g_{5}:=f_{12}$ $=-5 b d$.

## LEMMA 8.5.

For any finite input set $F_{0}=\left\{f_{1}, \ldots, f_{m}\right\}$ of polynomials, Algorithm 8.3 always terminates. PROOF. We consider the $\mathrm{F}_{\mathrm{k}}$ 's as multisets of polynomials, which are ordered by the multiset ordering << induced by the ordering < on polynomials ( see Definition 8.2 ). Since $<$ is well-founded, the multiset extension $\ll$ is also well-founded.
For the Cases 1, 2, 3 and 4.1.1, $\mathrm{F}_{\mathrm{k}} \gg \mathrm{F}_{\mathrm{k}+1}$. Case 4.1.2 and the according subcase of 4.1.3 can not occur infinitely often in successive steps, because then $g_{1}>g_{2}>g_{3}>\ldots$ would be an infinite descending <-chain. That means, that after finitely many steps $g_{i}=0$ or Case 4.1.1, the other subcase of 4.1.3 or Case 4.2 occur.
For the Cases 4.1.3 and 4.2, $\mathrm{F}_{\mathrm{k}+1}$ is larger than $\mathrm{F}_{\mathrm{k}}$. But these cases can only occur finitely often during the whole run of the algorithm. First note, that all words $t$ occurring in some polynomial of some $F_{k}$ satisfy $t \leq \max \left\{H W\left(f_{1}\right), \ldots, H W\left(f_{m}\right)\right\}$. Since $<$ is admissible, there are only finitely many words with this property. Hence Case 4.2 can only occur finitely often. Case 4.1.3 - where a head term, which has disappeared in some former step, appears again - can only occur finitely often for a certain term, because the absolut value of the head coefficient gets smaller each time.

Before we can state the next lemma, we have to introduce a new notation (or rather an abuse of the usual notation ). Let $F$ be a finite set of polynomials. The expression

$$
\mathrm{f}=\sum_{\mathrm{h}_{\mathrm{i}} \in \mathrm{~F}} \mathrm{~h}_{\mathrm{i}} \cdot \mathrm{a}_{\mathrm{i}},
$$

should be interpreted as follows: the $a_{i}$ are monomials in $\mathbb{Z}<X_{1}, \ldots, X_{n}>, f$ is a finite sum of the polynomials $h_{i} \cdot a_{i}$, but an element of $F$ may occur more than once in this sum and each occurrence may have a different coefficient $\mathrm{a}_{\mathrm{i}}$.

LEMMA 8.6. Let $t \in W_{n}$ be a word and $\mathrm{F}_{\mathrm{k}}$ be the set of polynomials obtained after some iterations of Algorithm 8.3. Assume that $h$ is a polynomial and that $h=\Sigma_{h_{i} \in F_{k}} h_{i} \cdot a_{i}$ for monomials $a_{i}$ with $H W\left(h_{i} \cdot a_{i}\right)<t$. Then $h=\Sigma_{h_{i}} \in F_{k+1} h_{i} \cdot b_{i}$ for monomials $b_{i}$ with $H W\left(h_{i}^{\prime} \cdot b_{i}\right)<t$.
PROOF. For the Cases 4.1.3 and 4.2 we have $\mathrm{F}_{\mathrm{k}} \subseteq \mathrm{F}_{\mathrm{k}+1}$ and thus we can use the given sum. In Case 1, $\mathrm{F}_{\mathrm{k}+1}:=\left(\mathrm{F}_{\mathrm{k}} \backslash\{\mathrm{f}\}\right) \cup\left\{\mathrm{f}_{1}\right\}$ and $\mathrm{f}=\mathrm{f}_{1}+\mathrm{g} \cdot \mathrm{c}$. In addition we have $\mathrm{g} \in \mathrm{F}_{\mathrm{k}}$ and $H W(g)=H W(f) \geq H W\left(f_{1}\right)$. Thus a term $f \cdot a_{j}$ in the sum $h=\Sigma_{h_{i} \in F_{k}} h_{i} \cdot a_{i}$ can be replaced by $f_{1} \cdot a_{j}+g \cdot c a_{j}$. The other cases can be treated similar.

LEMMA 8.7. Let $G$ be the output of Algorithm 8.3 ( i.e. the actual set $F_{k}$, when the algorithm terminates ) and let $f=a \cdot t+R(f)$ and $g=b \cdot t r+R(g)$ be elements of $G$.

Then the following holds:
(1) $a=b c$ for some $c \in \mathbb{Z},|c| \neq 1$ and $r \neq 1$.
(2) The S-polynomial $g_{1}:=f \cdot r-g \cdot c=R(f) \cdot r-R(g) \cdot c$ can be obtained as a finite sum

$$
g_{1}=\sum_{h_{i} \in G} h_{i} \cdot a_{i}
$$

where the $\mathrm{a}_{\mathrm{i}}$ are monomials in $\mathbb{Z}<\mathrm{X}_{1}, \ldots, \mathrm{X}_{\mathrm{n}}>$ and $H W\left(\mathrm{~h}_{\mathrm{i}} \cdot \mathrm{a}_{\mathrm{i}}\right) \leq \mathrm{HW}\left(\mathrm{g}_{1}\right)<\mathrm{HW}(\mathrm{g})=$ HW(f.r).
PROOF. Since Algorithm 8.3 has terminated, the index pair corresponding to $f$ and $g$ is marked. Thus for some $k, f$ and $g$ are in $F_{k}$ and they are selected by the algorithm.
(1) is satisfied, since only in Case 4 both $f$ and $g$ remain in $F_{k+1}$.
(2) In Case 4 we have $g_{1}:=\mathrm{f} \cdot \mathrm{r}-\mathrm{g} \cdot \mathrm{c}=\mathrm{R}(\mathrm{f}) \cdot \mathrm{r}-\mathrm{R}(\mathrm{g}) \cdot \mathrm{c}$ and thus $\mathrm{HW}\left(\mathrm{g}_{1}\right)<\mathrm{HW}(\mathrm{g})=$ $\mathrm{HW}(\mathrm{f} \cdot \mathrm{r})=\mathrm{tr}$. There is some $\mathrm{g}_{\mathrm{i}}$ such that $\left.\mathrm{g}_{1}{ }^{*}\right)_{\mathrm{F}_{k}} \mathrm{~g}_{\mathrm{i}} \quad$ ( see Case 4.1.2 and the first subcase of 4.1.3 ) and $g_{i} \in F_{k+1}$ or $g_{i}=0$. Hence $H W\left(g_{i}\right) \leq H W\left(g_{1}\right)$ and $g_{1}=g_{i}+\Sigma_{h_{i} \in F_{k}} h_{i} \cdot a_{i}$ for monomials $a_{i}$ with $H W\left(h_{i} \cdot a_{i}\right) \leq H W\left(g_{1}\right)$. Lemma 8.6 yields $g_{1}=g_{i}+\Sigma_{h_{i} \in F_{k+1}} h_{i} \cdot b_{i}$ for monomials $b_{i}$ with $H W\left(h_{i} \cdot b_{i}\right) \leq H W\left(g_{1}\right)$ and since $g_{i} \in F_{k+1}$ or $g_{i}=0$ we have $g_{1}=$ $\Sigma_{h_{i}{ }^{n} \in F_{k+1}} h_{i} " \cdot c_{i}$ for monomials $c_{i}$ with HW $\left(h_{i} " \cdot c_{i}\right) \leq H W\left(g_{1}\right)$. By Lemma 8.6, $g_{1}$ can be represented by such a sum for all $\mathrm{F}_{\mathrm{m}}$ with $\mathrm{m} \geq \mathrm{k}+1$. Thus we have proved the lemma.

Let $\mathrm{F} \subseteq \mathbb{Z}<\mathrm{X}_{1}, \ldots, \mathrm{X}_{\mathrm{n}}>$ be a set of polynomials. In the following $<\mathrm{F}>$ denotes the right ideal generated by F

## LEMMA 8.8.

Let $G$ be the output of Algorithm 8.3 if started with input $F_{0}$. Then $\langle G\rangle=\left\langle F_{0}\right\rangle$.
PROOF. It is easy to see that, for any $k,\left\langle\mathrm{~F}_{\mathrm{k}}\right\rangle=\left\langle\mathrm{F}_{\mathrm{k}+1}\right\rangle$.

This lemma and the next proposition shows, that it is reasonable to call the result of Algorithm 8.3 a Gröbner base.

PROPOSITION 8.9. Let $G$ be the output of Algorithm 8.3. Then any $f \in<G>$ can be reduced to 0 w.r.t. $\rightarrow_{G}$.
PROOF. The proof is similar to the proof of Lemma 2.4 in Mora (1985).
Obviously, $\mathrm{f} \in\langle\mathrm{G}\rangle$ means $\mathrm{f}=\Sigma_{\mathrm{g} i \in \mathrm{G}} \mathrm{g}_{\mathrm{i}} \cdot \mathrm{a}_{\mathrm{i}}$ for some monomials $\mathrm{a}_{\mathrm{i}}$. Let $\mathrm{t}:=\max \{\ldots$ $\left.\operatorname{HW}\left(g_{i} \cdot a_{i}\right) \ldots\right\}$ and $I:=\left\{i ; H W\left(g_{i} \cdot a_{i}\right)=t\right\}$.
Case 1. $|\mathrm{II}|=1$.
Then $H W(f)=t$ and $\left(\right.$ for $I=\{j\}$ and $\left.a_{j}=c_{j} \cdot r_{j}\left(c_{j} \in \mathbb{Z}, r_{j} \in W_{n}\right)\right) H W(f)=t=$ $\mathrm{HW}\left(\mathrm{g}_{\mathrm{j}}\right) \cdot \mathrm{r}_{\mathrm{j}}$ and $\mathrm{HC}(\mathrm{f})=\mathrm{HC}\left(\mathrm{g}_{\mathrm{j}}\right) \cdot \mathrm{c}_{\mathrm{j}}$. Hence f can be reduced by $\mathrm{g}_{\mathrm{j}}$ to the smaller polynomial $f_{1}:=f-g_{j} \cdot a_{j} \in\langle G\rangle$. By Induction we get $f_{1}{ }_{\mathrm{H}}{ }_{G} 0$ and thus $f \rightarrow_{G} f_{i}{ }_{G}^{*} 0$.

Case 2. |II > 1.
Let $i, j$ be two different elements of $I$ and let $a_{i}=c_{i} \cdot r_{i}, a_{j}=c_{j} \cdot r_{j}\left(c_{i}, c_{j} \in \mathbb{Z}, r_{i}, r_{j} \in W_{n}\right)$. Since $\mathrm{HW}\left(\mathrm{g}_{\mathrm{i}}\right) \cdot \mathrm{r}_{\mathrm{i}}=\mathrm{t}=\mathrm{HW}\left(\mathrm{g}_{\mathrm{j}}\right) \cdot \mathrm{r}_{\mathrm{j}}$, either $\mathrm{HW}\left(\mathrm{g}_{\mathrm{i}}\right)$ is a prefix of $\mathrm{HW}\left(\mathrm{g}_{\mathrm{j}}\right)$ or vice versa. Without loss of generality we assume $\mathrm{HW}\left(\mathrm{g}_{\mathrm{i}}\right)=\mathrm{HW}\left(\mathrm{g}_{\mathrm{j}}\right) \cdot \mathrm{r}$ for some $\mathrm{r} \in \mathrm{W}_{\mathrm{n}}$. By Lemma 8.7, $H C\left(g_{j}\right)=H C\left(g_{i}\right) \cdot c$ for some $c \in \mathbb{Z}$ and $g_{j} \cdot r-g_{i} \cdot c=\Sigma_{h_{k} \in G} h_{k} \cdot b_{k}$, where $H W\left(h_{k} \cdot b_{k}\right)<$ $H W\left(g_{i}\right)=H W\left(g_{j} \cdot r\right)$. Hence $g_{j} \cdot r_{j}-g_{i} \cdot r_{i} c=\left(g_{j} \cdot r-g_{i} \cdot c\right) \cdot r_{i}=\Sigma_{h_{k} \in G} h_{k} \cdot\left(b_{k} r_{i}\right)$, where $\operatorname{HW}\left(h_{k} \cdot\left(b_{k} r_{i}\right)\right)<H W\left(g_{i}\right) \cdot r_{i}=t$.
Now $f=\left(g_{j} \cdot r_{j}-g_{i} \cdot r_{i} c\right) \cdot c_{j}+g_{i}\left(c_{i}+c c_{j}\right) r_{i}+\Sigma_{m \neq i, j} g_{m} \cdot a_{m}$
$=\Sigma_{h_{k} \in G} h_{k} \cdot\left(b_{k} c_{j} r_{i}\right)+g_{i}\left(c_{i}+c c_{j}\right) r_{i}+\Sigma_{m \neq i, j} g_{m} \cdot a_{m}$ yields a representation of $f$ as a sum, where III is smaller.

COROLLARY 8.10. The membership problem for finitely generated right ideals in $\mathbb{Z}<X_{1}, \ldots, X_{n}>$ is decidable.
PROOF. Let $I=<\left\{f_{1}, \ldots, f_{m}\right\}>$ be a finitely generated right ideal in $\left.\mathbb{Z}<X_{1}, \ldots, X_{n}\right\rangle$. We apply Algorithm 8.3 to $F_{0}=\left\{f_{1}, \ldots, f_{m}\right\}$ and get a set $G$ of polynomials. Now $f \in I$ iff $f$ can be reduced to 0 w.r.t. $\rightarrow_{G}$. If $f$ is $\rightarrow_{G}$-irreducible, then $f \in I$ iff $f=0$. Otherwise we can effectively find some $g$ such that $f \rightarrow_{G} g$ and $f \in I$ iff $g \in I$. Thus Corollary 8.10 is proved by induction.

## 9. Solving Linear Equations in $\mathbb{Z}<X_{1}, \ldots, X_{n}>$

In the previous section we have shown, how to compute "Gröbner bases" for finitely generated right ideals in $\mathbb{Z}<X_{1}, \ldots, X_{n}>$. In this section these bases are used to solve linear equations in $\mathbb{Z}<X_{1}, \ldots, X_{n}>$. The method is very similar to that described in Section 6.
Let (*) $f_{1} x_{1}+\ldots+f_{r} x_{r}=f_{0}$ be an (inhomogeneous ) linear equation in $\mathbb{Z}<X_{1}, \ldots, X_{n}>$. We have to find one solution for ( $*$ ) and finitely many generators of the solutions of the homogeneous equation $(* *) \mathrm{f}_{1} \mathrm{x}_{1}+\ldots+\mathrm{f}_{\mathrm{r}} \mathrm{x}_{\mathrm{r}}=0$.
Let $G=\left\{g_{1}, \ldots g_{s}\right\}$ be the output of Algorithm 8.3. when started with input $\left\{f_{1}, \ldots, f_{r}\right\}$. There exist an r×s-matrix $P$ and an $s \times r$-matrix $Q$ with entries in $\mathbb{Z}<X_{1}, \ldots, X_{n}>$ such that $\underline{\mathrm{f}} \cdot \mathrm{P}=\mathrm{g}$ and $\mathrm{g} \cdot \mathrm{Q}=\underline{\mathrm{f}}$. This matrices can be obtained as by-products of Algorithm 8.3.
Obviously, (*) has a solution iff $\mathrm{f}_{0} \in\left\langle\left\{\mathrm{f}_{1}, \ldots, \mathrm{f}_{\mathrm{r}}\right\}\right\rangle=\langle\mathrm{G}\rangle$. Hence, if (*) has a solution, Proposition 8.9 implies that $f_{0}$ reduces to 0 w.r.t $\rightarrow_{G}$. This yields $p_{1}, \ldots, p_{s} \in \mathbb{Z}\left[X_{1}, \ldots, X_{n}\right]$ such that $\mathrm{g} \cdot \mathrm{p}=\mathrm{f}_{0}$. But then $\mathrm{P} \cdot \mathrm{lp}$ is a solution of $(*)$.
We now assume that we already have finitely many generators $\mathrm{lz}{ }^{(1)}, \ldots, z^{(L)}$ of the set of solutions of the equation

$$
(++) g_{1} x_{1}+\ldots+g_{s} x_{s}=0
$$

As in Section 6 one can show

LEMMA 9.1. The vectors $P \cdot\left|z^{(1)}, \ldots, P \cdot\right| z^{(L)}$ and the columns of the matrix $P Q-E_{r}$ are solutions of (**) and they generate all solutions of this equation.

We now show how to compute the finitely many generators of the solutions of $(++)$. If there do not exist $i, j(i \neq j)$ such that $H W\left(g_{i}\right)=H W\left(g_{j}\right) \cdot r$ for some $r \in W_{n}$, the equation $(++)$ has no nontrivial solutions. Otherwise, let $i, j(i \neq j)$ be indices, such that $H W\left(g_{i}\right)=H W\left(g_{j}\right) \cdot r$ for some $r \in W_{n}$.
By Lemma 8.7, $\mathrm{HC}\left(\mathrm{g}_{\mathrm{j}}\right)=\mathrm{HC}\left(\mathrm{g}_{\mathrm{i}}\right) \cdot \mathrm{c}$ for some $\mathrm{c} \in \mathbb{Z}, \mathrm{r} \neq 1$ and

$$
g_{j} \cdot \mathrm{r}-\mathrm{g}_{\mathrm{i}} \cdot \mathrm{c}=\sum_{\mathrm{k}=1}^{\mathrm{k}=\mathrm{r}} \mathrm{~g}_{\mathrm{k}} \cdot \mathrm{~h}_{\mathrm{k}}
$$

for polynomials $h_{k} \in \mathbb{Z}<X_{1}, \ldots, X_{n}>$ with $H W\left(g_{k} \cdot h_{k}\right)<H W\left(g_{i}\right)$. Obviously, $h_{i}$ has to be 0 .
If we define $q_{k}:=h_{k}$ for $k \neq i, j, q_{i}:=h_{i}+c=c$ and $q_{j}:=h_{j}-r$, then $q_{i j}:=\left(q_{1}, \ldots, q_{S}\right)^{T}$ is a solution of $(++)$.

LEMMA 9.2. The finitely many vectors $\mathrm{lq}_{\mathrm{ij}}$ generate all solutions of (++).
PROOF. Let $\mid \mathrm{p}=\left(\mathrm{p}_{1}, \ldots, \mathrm{p}_{\mathrm{s}}\right)^{\mathrm{T}}$ be a nontrivial solution of $(++)$. The complexity of such a solution is given by $(t, \alpha)$, where $t:=\max \left\{H W\left(g_{i} p_{i}\right) ; 1 \leq i \leq s\right\}$ and $\alpha:=\mid\{i ; 1 \leq i \leq s$ and $\left.\mathrm{HW}\left(\mathrm{g}_{\mathrm{i}} \mathrm{p}_{\mathrm{i}}\right)=\mathrm{t}\right\}$.
Since $\mathrm{g} \cdot \mathrm{p}=0$ and lp is not trivial, $\alpha$ has to be greater than 1 . Hence there exist $\mathrm{i}, \mathrm{j}(\mathrm{i} \neq \mathrm{j})$ such that $\mathrm{HW}\left(\mathrm{g}_{\mathrm{i}}\right) \mathrm{HW}\left(\mathrm{p}_{\mathrm{i}}\right)=\mathrm{t}=\mathrm{HW}\left(\mathrm{g}_{\mathrm{j}}\right) \mathrm{HW}\left(\mathrm{p}_{\mathrm{j}}\right)$. Without loss of generality we assume that $\mathrm{HW}\left(\mathrm{g}_{\mathrm{j}}\right)$ is a prefix of $\mathrm{HW}\left(\mathrm{g}_{\mathrm{i}}\right)$. Thus $\mathrm{HW}\left(\mathrm{g}_{\mathrm{i}}\right)=\mathrm{HW}\left(\mathrm{g}_{\mathrm{j}}\right) \cdot \mathrm{r}$ and $\mathrm{HC}\left(\mathrm{g}_{\mathrm{j}}\right)=\mathrm{HC}\left(\mathrm{g}_{\mathrm{i}}\right) \cdot \mathrm{c}$ for some $r \in W_{n}$ and $c \in \mathbb{Z}$ and $H W\left(p_{j}\right)=r \cdot H W\left(p_{i}\right)$. Let $c_{i}:=H C\left(p_{i}\right)$ and $c_{j}:=H C\left(p_{j}\right)$.
The vector $\mid \mathrm{q}_{\mathrm{ij}}$, which was defined above, is a solution of $(++)$. We define a new solution $\left(p_{1}{ }^{\prime}, \ldots, p_{s}{ }^{\prime}\right)^{T}=\left|p^{\prime}:=|p+| q_{i j} \cdot c_{j} H W\left(p_{i}\right)\right.$ and show that it has smaller complexity than $| \mathrm{p}$. To that purpose we have to consider the words $\mathrm{HW}\left(\mathrm{g}_{\mathrm{k}} \mathrm{p}_{\mathrm{k}}{ }^{\prime}\right)$ for all $\mathrm{k}, 1 \leq \mathrm{k} \leq \mathrm{s}$.
CASE 1. $\mathrm{k} \neq \mathrm{i}, \mathrm{j}$.
We have $g_{k} p_{k}^{\prime}=g_{k} p_{k}+g_{k} h_{k} c \cdot H W\left(p_{i}\right)$ and $H W\left(g_{k} \cdot h_{k}\right)<H W\left(g_{i}\right)$. This implies that $\operatorname{HW}\left(\mathrm{g}_{\mathrm{k}} \mathrm{h}_{\mathrm{k}} \mathrm{c}_{\mathrm{j}} \mathrm{HW}\left(\mathrm{p}_{\mathrm{i}}\right)\right)<\operatorname{HW}\left(\mathrm{g}_{\mathrm{i}}\right) \mathrm{HW}\left(\mathrm{p}_{\mathrm{i}}\right)=\mathrm{t}$. Thus $\mathrm{HW}\left(\mathrm{g}_{\mathrm{k}} \mathrm{p}_{\mathrm{k}}{ }^{\prime}\right)=\mathrm{t}$, if $\mathrm{HW}\left(\mathrm{g}_{\mathrm{k}} \mathrm{p}_{\mathrm{k}}\right)=\mathrm{t}$, and otherwise, $\mathrm{HW}\left(\mathrm{g}_{\mathrm{k}} \mathrm{p}_{\mathbf{k}}\right)<\mathrm{t}$.
Case 2. $\mathrm{k}=\mathrm{i}$.
We have $\mathrm{g}_{\mathrm{i}} \mathrm{p}_{\mathrm{i}}{ }^{\prime}=\mathrm{g}_{\mathrm{i}} \mathrm{p}_{\mathrm{i}}+\mathrm{g}_{\mathrm{i}} \mathrm{c} \mathrm{c}_{\mathrm{j}} \mathrm{HW}\left(\mathrm{p}_{\mathrm{i}}\right)$. Hence $\mathrm{HW}\left(\mathrm{g}_{\mathrm{i}} \mathrm{p}_{\mathrm{i}}{ }^{\prime}\right)=\mathrm{t}$ if $\mathrm{c}_{\mathrm{i}}+\mathrm{cc}_{\mathrm{j}} \neq 0$ and $\mathrm{HW}\left(\mathrm{g}_{\mathrm{i}} \mathrm{p}_{\mathrm{i}}{ }^{\prime}\right)<\mathrm{t}$ if $\mathrm{c}_{\mathrm{i}}+\mathrm{cc}_{\mathrm{j}}=0$.
Case 3. $\mathrm{k}=\mathrm{j}$.
$g_{j} p_{j}^{\prime}=g_{j} p_{j}+g_{j} h_{j} c_{j} H W\left(p_{i}\right)-g_{j}{ }_{j}{ }_{j} H W\left(p_{i}\right)$

$$
\begin{aligned}
& =H C\left(g_{j}\right) c_{j} t+R\left(g_{j} p_{j}\right)+g_{j} h_{j} c_{j} H W\left(p_{i}\right)-H C\left(g_{j}\right) c_{j} H W\left(g_{j}\right) r H W\left(p_{i}\right)-R\left(g_{j}\right) r c_{j} H W\left(p_{i}\right) \\
& =R\left(g_{j} p_{j}\right)+g_{j} h_{j} c_{j} H W\left(p_{i}\right)-R\left(g_{j}\right) r c_{j} H W\left(p_{i}\right), \text { since } \operatorname{rHW}\left(p_{i}\right)=H W\left(g_{j}\right) .
\end{aligned}
$$

This shows that $\mathrm{HW}\left(\mathrm{g}_{\mathrm{j}} \mathrm{p}_{\mathrm{j}}{ }^{\prime}\right)<\mathrm{t}$.

Thus we have seen the t the complexity of the solution $\mid \mathrm{p}$ ' is smaller then the complexity of lp and the lemma is proved by induction.

EXAMPLE 9.3. As an example we consider the homogeneous linear equation $f_{1} x_{1}+\ldots+$ $\mathrm{f}_{4} \mathrm{x}_{4}=0$ in $\mathbb{Z}<a, b, c, d>$ for the polynomials $\mathrm{f}_{1}=2 \mathrm{abc}-\mathrm{bc}, \mathrm{f}_{2}=3 \mathrm{ab}-2 \mathrm{~b}, \mathrm{f}_{3}=5 \mathrm{abd}-\mathrm{bc}$ and $\mathrm{f}_{4}=\mathrm{bc}-5 \mathrm{bd}$ of Example 8.4.
We have seen that Algorithm 8.3 terminates with $G=\left\{g_{1}, g_{2}, g_{3}, g_{4}, g_{5}\right\}$, where $g_{1}=$ $3 a b-2 b, g_{2}=a b c-b c, g_{3}=b c, g_{4}=a b d-4 b d+b c$ and $g_{5}=-5 b d$. The transformation matrices $P, Q$ such that $\underline{f} \cdot P=g$ and $g \cdot Q=\underline{f}$ are

$$
\mathrm{Q}=\left(\begin{array}{cccc}
0 & 1 & \mathrm{~d} & 0 \\
2 & 0 & 0 & 0 \\
1 & 0 & -3 & 1 \\
0 & 0 & 2 & 0 \\
0 & 0 & -2 & 1
\end{array}\right) \quad \text { and } \quad P=\left(\begin{array}{ccccc}
0 & -1 & 3 & 0 & 0 \\
1 & \mathrm{c} & -2 \mathrm{c} & 2 \mathrm{~d} & -5 \mathrm{~d} \\
0 & 0 & 0 & -1 & 3 \\
0 & 0 & 0 & 0 & 3
\end{array}\right)
$$

All solutions of the equation $g_{1} x_{1}+\ldots+g_{5} x_{5}=0$ are generated by $\mathrm{lq}_{1,2}$ and $\mathrm{lq}_{1,4}$ :
(1) $g_{1} \cdot c-g_{2} \cdot 3=g_{3}$ and thus $\operatorname{lq}_{1,2}=(-c, 3,1,0,0)^{T}$.
(2) $g_{1} \cdot d-g_{4} \cdot 3=f_{11}=f_{12}-f_{4} \cdot 3=f_{12}-\left(f_{6}+f_{12}\right) \cdot 3=f_{12} \cdot(-2)+f_{6}(-3)=g_{5} \cdot(-2)+g_{3}(-3)$ and thus $\mathrm{Iq}_{1,4}=(-\mathrm{d}, 0,-3,3,-2)^{\mathrm{T}}$.
We now apply $P$, to get the corresponding solutions of $f_{1} x_{1}+\ldots+f_{4} x_{4}=0$ :
$P \cdot \mid q_{1,2}=(0,0,0,0)^{T}$ and $P \cdot \mid q_{1,4}=(-9,6 c+15 d,-9,-6)^{T}$.
The matrix $\mathrm{PQ}-\mathrm{E}_{4}$ is $\left(\begin{array}{cccc}0 & 0 & -9 & 3 \\ 0 & 0 & 6 \mathrm{c}+15 \mathrm{~d} & -2 \mathrm{c}-5 \mathrm{~d} \\ 0 & 0 & -9 & 3 \\ 0 & 0 & -6 & 2\end{array}\right)$.
This yields the new solution ( $3,-2 \mathrm{c}-5 \mathrm{~d}, 3,2)^{\mathrm{T}}$ and since $\mathrm{lq}_{1,4}=(3,-2 \mathrm{c}-5 \mathrm{~d}, 3,2)^{\mathrm{T}} \cdot(-3)$, the solution $(3,-2 c-5 d, 3,2)^{T}$ generates all solutions of $f_{1} x_{1}+\ldots+f_{4} x_{4}=0$ in $\mathbb{Z}<\mathrm{a}, \mathrm{b}, \mathrm{c}, \mathrm{d}>$.

## 10. Conclusion

The categorical reformulation of E-unification allows to characterize the class of commutative theories by properties of the category $\mathrm{C}(\mathrm{E})$ of finitely generated E -free objects: $C(E)$ has to be a semiadditive category. The definition of semiadditive categories provides an algebraic structure on the morphism sets, which can be used to obtain algebraic characterizations of the unification types. This shows the connection between unification in commutative theories and equation solving in linear algebra. The very common syntactic approach to equational unification, which only uses the defining axioms, is thus replaced by a more semantic approach, which works with algebraic properties of the defined algebras.

Hence unification algorithms for the commutative theory AGnHC, i.e. the theory of abelian groups with n commuting homomorphisms, can be derived with the help of wellknown algebraic methods ( e.g. Gröbner Base algorithms) to solve linear equations in $\mathbb{Z}\left[\mathrm{X}_{1}, \ldots, \mathrm{X}_{\mathrm{n}}\right]$. In order to obtain a unification algorithm for the theory AGnH of abelian groups with n non-commuting homomorphisms, we developed a Gröbner base algorithm for the ring $\mathbb{Z}<X_{1}, \ldots, X_{n}>$ of polynomials over $\mathbb{Z}$ in $n$ non-commuting indeterminates. Since Dicksons Lemma ( Dickson (1913)), which is used for $\mathbb{Z}\left[X_{1}, \ldots, X_{n}\right]$ to prove termination of the Gröbner Base algorithm, does not hold for $\mathbb{Z}<X_{1}, \ldots, X_{n}>$, we had to be very careful to obtain a terminating algorithm. As in the commutative case, the performance of the algorithm depends on the choice of the admissible ordering. Hence it would be interesting to have a complete characterization of all admissible orderings for $\mathbb{Z}\left\langle X_{1}, \ldots, X_{n}>\right.$.

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[^0]:    1) This research was done while the author was still at the IMMD 1, University Eriangen.
