# ROSENBERG'S CHARACTERIZATION OF MAXIMAL CLONES 

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#### Abstract

We will give a proof of I. G. Rosenberg's characterization of maximal clones, first published in [11]. The theorem lists six types of relations on a finite set such that a clone over this set is maximal if and only if it contains just the functions preserving one of the relations of the list. In Universal Algebra, this translates immediately into a characterization of the finite preprimal algebras: A finite algebra is preprimal if and only if its term operations are exactly the functions preserving a relation of one of the six types listed in the theorem. The difficult part of the proof is to show that all maximal clones or preprimal algebras respectively are of that form. This follows from, and, as we will also demonstrate, is indeed equivalent to, a characterization of primal algebras: We will show that the primal algebras are exactly those whose term operations do not preserve any of the relations on the list.


## Preface

A clone (closed operation network) $\mathcal{C}$ over a set $A$ is a set of operations on this set which contains the projections and which is closed under compositions. The set of all clones over $A$ forms a lattice Clone $(A)$ with respect to inclusion, and a clone is called maximal if and only if it is a dual atom in $\operatorname{Clone}(A)$.

It is a fact that if $A$ is finite, then every clone is contained in a maximal clone and the maximal clones are finite in number. In his work [11] I. G. Rosenberg gave a characterization of the maximal clones over a a finite base set in terms of relations: The theorem lists six types of relations on $A$ such that a clone is maximal if and only if it is just the set of functions preserving one of the relations of the list.

However, the original proof of this deep theorem is quite technical and hard to follow. It is the aim of the present work to provide a shorter and somewhat more understandable proof.

Our proof is based on the one by R. W. Quackenbush in [9], who showed the more difficult implication of the theorem, namely that every maximal clone is of the form described before. It draws heavily on results of R. W. Quackenbush [10] on algebras with minimal spectrum, of H. P. Gumm [8] on algebras in permutable varieties, and of A. Foster and A. Pixley [6] on primality. Also a part of the original proof of I. G. Rosenberg has been included. We would like to add that there exists another new proof of the difficult implication of the theorem by V .
A. Buevich in [2].

This thesis has been divided into three chapters. In the first chapter, we introduce the theorem and explain the connection between maximal clones and preprimal algebras. Chapter 2 contains the proof of half of the equivalence: Every maximal clone is a set of functions preserving one of the relations listed in the theorem. Chapter 3 is devoted to the proof of the converse statement that all relations of the list yield a maximal clone.

All global conventions regarding notation will be made in the first chapter together with the basic definitions, and additional conventions will be introduced in Notations 2.0.10, 2.1.5 and 3.0.14. We tried to keep this work self-contained, the reader is assumed to be familiar only with the rudiments of Universal Algebra, lattice theory, and some basic facts about groups and fields; information on clones can be found in [12].

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## Chapter 1

## Rosenberg's preprimal algebra characterization

We will state the characterization of the maximal clones and provide the reader with the necessary definitions. Moreover, the connection between another possible viewpoint of the theorem, namely the characterization of finite preprimal algebras, and the theorem itself as a statement about clones will be explained.
1.0.1 Definition. Let $A$ be a set and denote by $\mathcal{F}_{n}$ the set of all $n$-ary functions on $A$. Then $\mathcal{F}=\bigcup_{n=0}^{\infty} \mathcal{F}_{n}$ is the set of all functions on $A$ of arbitrary arity. A clone is a subset of $\mathcal{F}$ which is closed under compositions and which contains all projections. The set of all clones on $A$ form a lattice $C l o n e(A)$ with respect to inclusion. A clone is called maximal iff it is maximal in $\operatorname{Clone}(A) \backslash\{\mathcal{F}\}$.

In order to bring these definitions into the context of Universal Algebra, one can think of a clone $\mathcal{C}$ on $A$ as the set of term operations of the algebra $\mathfrak{A}=(A, \mathcal{C})$. Conversely, given an algebra $\mathfrak{A}=(A, F)$, the term operations $\mathfrak{T}(F)$ form a clone over $A$. This interpretation of clones makes sense, for it provides the possibility of making use of the existing apparatus of Universal Algebra, e.g. congruence relations. It is for this reason that we will talk about algebras rather than about clones for the biggest part of our proof.
1.0.2 Definition. An algebra $\mathfrak{A}$ is primal iff every function on $A$ is a term operation of $\mathfrak{A} ; \mathfrak{A}$ is preprimal iff it is not primal but for any function $f$ not a term operation of $\mathfrak{A},(A, F \cup\{f\})$ is primal.

By the previous discussion, maximal clones correspond to preprimal algebras and vice-versa. Let $\mathcal{R}_{n}$ be the set of all $n$-ary relations on $A$; then $\mathcal{R}=\bigcup_{n=1}^{\infty} \mathcal{R}_{n}$ is the set of all relations on $A$ of arbitrary finite arity. We define for an arbitrary set $R \subseteq \mathcal{R}$ of relations on $A$ the set of polymorphisms $\operatorname{Pol}(R)$, that is, if we write $R_{k}$ for the $k$-ary relations in $R$ and $a_{1}, \ldots, a_{n}$ for
the coordinates of an $n$-tuple $a$,

$$
\begin{aligned}
& \operatorname{Pol}(R)=\bigcup_{n=0}^{\infty}\left\{f \in \mathcal{F}_{n}: \forall k \geq 0 \forall \rho \in R_{k} \forall r_{1}, \ldots, r_{n} \in \rho\right. \\
&\left.\left(\left(f\left(r_{11}, \ldots, r_{n 1}\right), \ldots, f\left(r_{1 k}, \ldots, r_{n k}\right)\right) \in \rho\right)\right\} .
\end{aligned}
$$

With this definition, Rosenberg's theorem states that a clone over a finite set $A$ is maximal iff it is of the form $\operatorname{Pol}(\{\rho\})$, where $\rho$ is a relation in one of six classes to be specified later. To formulate Rosenberg's theorem in detail, we need a couple of definitions.

For a function $f$ on $A$ define the graph of $f$ to be the set $\{(a, f(a)): a \in A\}$. Sometimes we will talk about a function and mean the graph of the function as a subset of $A^{2}$; confusion is unlikely since things should be clear from context.
A permutation $\pi$ is prime iff all cycles of $\pi$ have the same prime length.
We call a subset $\rho \subseteq A^{4}$ affine iff there is a binary operation + on $A$ such that $(A,+)$ is an abelian group and $(a, b, c, d) \in \rho \leftrightarrow a+b=c+d$ holds. An affine $\rho$ is prime iff $(A,+)$ is an abelian $p$-group for some prime $p$, that is, all elements of the group have the same prime order $p$.
For $h \geq 1$ a subset $\rho \subseteq A^{h}$ is totally symmetric iff for all permutations $\pi$ of $\{1, \ldots, h\}$ and all tuples $\left(a_{1}, \ldots, a_{h}\right) \in A^{h},\left(a_{1}, \ldots, a_{h}\right) \in \rho$ iff $\left(a_{\pi(1)}, \ldots, a_{\pi(n)}\right) \in \rho$. Define $\iota_{h}^{A} \subseteq A^{h}$ by

$$
\iota_{h}^{A}=\left\{\left(a_{1}, \ldots, a_{h}\right) \mid \exists i \exists j\left(i \neq j \wedge a_{i}=a_{j}\right)\right\} .
$$

Then $\rho$ is called totally reflexive iff $\iota_{h}^{A} \subseteq \rho$. Note that for $h=2$, totally reflexive means reflexive and totally symmetric means symmetric. If $\rho$ is totally reflexive and totally symmetric we define the center of $\rho$ to be the set

$$
C(\rho)=\left\{a \in A \mid \forall a_{2}, \ldots, a_{h} \in A\left(a, a_{2}, \ldots, a_{h}\right) \in \rho\right\} .
$$

We say that $\rho \subseteq A^{h}$ is central iff it is totally reflexive, totally symmetric and has a nonvoid center which is a proper subset of $A$. Note that $h \leq|A|$ as otherwise we would have $\rho \supseteq \iota_{h}^{A}=A^{h}$ and the center of $\rho$ would be trivial.
For an arbitrary set $S$ and $1 \leq r \leq \lambda$, denote the $r-$ th projection from $S^{\lambda}$ onto $S$ by $\pi_{r}^{\lambda}$. Now let $h=\{0,1, \ldots, h-1\}$ and define $\omega_{\lambda}$ to be the $h$-ary relation on $h^{\lambda}$ satisfying $\left(a_{1}, \ldots, a_{h}\right) \in \omega_{\lambda}$ iff for all $1 \leq r \leq \lambda,\left(\pi_{r}^{\lambda}\left(a_{1}\right), \ldots, \pi_{r}^{\lambda}\left(a_{h}\right)\right) \in \iota_{h}^{h}$. For $3 \leq h \leq|A|$, we call a $h$-ary relation $\rho$ on $A h$-regularly generated iff there exists a $\lambda \geq 1$ and a surjection $\varphi: A \rightarrow h^{\lambda}$ such that $\rho=\varphi^{-1}\left(\omega_{\lambda}\right)$. Note that for any relation, $h$-regularly generated implies totally reflexive and totally symmetric.

Now here comes the theorem.
1.0.3 Theorem (I. G. Rosenberg [11]). Let $1<|A|<\aleph_{0}$. A clone $\mathcal{C}$ on $A$ is maximal if and only if it is of the form $\operatorname{Pol}(\rho)$, where $\rho$ is an $h$-ary relation belonging to one of the following classes:

1. The set of all partial orders with least and greatest element
2. The set of all prime permutations
3. The set of all non-trivial equivalence relations
4. The set of all prime-affine relations
5. The set of all central relations
6. The set of all h-regularly generated relations

We will refer to the six classes as Rosenberg's list ( $R B L$ ) from now on. Then in the terminology of algebras, the theorem sounds like this.
1.0.4 Corollary. A finite non-trivial algebra $\mathfrak{A}$ is preprimal iff there exists a relation $\rho$ in $R B L$ such that $\mathfrak{T}(\mathfrak{A})=\operatorname{Pol}(\rho)$.
1.0.5 Remark. As with Rosenberg's theorem the maximal clones over a set $A$ with finite cardinality $\kappa$ are known, one can calculate their number $\eta_{\kappa}$. That number grows fast with the size $\kappa$. Here are values for a couple of cardinalities $\kappa$ :

| $\kappa$ | 2 | 3 | 4 | 5 | 6 | 7 |
| :---: | :---: | :---: | :---: | :---: | :---: | :---: |
| $\eta_{\kappa}$ | 5 | 18 | 82 | 643 | 15182 | 7848984 |

1.0.6 Remark. The clone lattice $\operatorname{Clone}(A)$ is countable only for $|A|=2$. For $|A| \geq 3$ we have already $|\operatorname{Clone}(A)|=2^{\aleph_{0}}$.

## Chapter 2

## Primal algebra characterization

We will prove the more difficult part of the equivalence by proving the following theorem.
2.0.7 Theorem. If a finite non-trivial algebra $\mathfrak{A}$ has no subalgebra of a finite power of $\mathfrak{A}$ belonging to $R B L$, then $\mathfrak{A}$ is primal.

The required implication in Rosenberg's theorem follows indeed.
2.0.8 Corollary. If a finite non-trivial algebra $\mathfrak{A}$ is preprimal then the set of term operations of $\mathfrak{A}$ is of the form $\operatorname{Pol}(\rho)$, where $\rho$ is a relation in $R B L$.

Proof. Since $\mathfrak{A}$ is not primal, by the last theorem there exists a subalgebra of a finite power of $\mathfrak{A}$ with universe $\rho$ in $R B L$; hence, the term operations satisfy $\mathfrak{T}(\mathfrak{A}) \subseteq \operatorname{Pol}(\rho)$. But as $\operatorname{Pol}(\rho)$ is closed under composition and projections and as $\mathfrak{A}$ is preprimal, $\mathfrak{T}(\mathfrak{A})=\operatorname{Pol}(\rho)$.

The corollary is in fact equivalent to the theorem.
2.0.9 Theorem. If all finite non-trivial preprimal algebras $\mathfrak{A}$ satisfy $\mathfrak{T}(\mathfrak{A})=\operatorname{Pol}(\rho)$, where $\rho$ is a relation in $R B L$, then every finite non-trivial algebra which preserves no relation belonging to $R B L$ is primal.

Proof. Let $\mathfrak{A}$ be a finite non-trivial algebra preserving no relation belonging to $R B L$. Then the clone $\mathfrak{T}(\mathfrak{A})$ is contained in no clone of the form $\operatorname{Pol}(\rho), \rho \in R B L$. But since all maximal clones are of that form and since Clone $(A)$ is dually atomic (see [12]), this means that $\mathfrak{T}(\mathfrak{A})$ must be the greatest element in that lattice and thus the clone of all functions on $A$. Hence, $\mathfrak{A}$ is primal.

To prove Theorem 2.0.7, we will first show that the hypotheses imply that all subalgebras of finite powers of $\mathfrak{A}$ have cardinality a power of the cardinality of $\mathfrak{A}$, which is a result by $R$. Quackenbush in [9]. R. W. Quackenbush also essentially showed in [10] that then the algebra generates a congruence permutable variety; we will follow his proof in the beginning, but then
use a slightly different approach to prove this, combining works of D. Clark and P. Krauss in [5] and of I. Chajda and G. Eigenthaler in [4]. Following H. P. Gumm in [8] and then H. Werner in [13] we will conclude that all powers of $\mathfrak{A}$ can only have factor congruences, which trivially implies that the equational class generated by $\mathfrak{A}$ is congruence distributive. A criterion for primality due to A. Foster and A. Pixley [6] will finally conclude the proof. Here is a summary of which implications we will prove; it might be helpful to look at it from time to time. The notions which occur in those implications will be defined in the respective sections.

- If $\mathfrak{A}$ is a finite non-trivial algebra having no subalgebra of a power of $\mathfrak{A}$ belonging to $R B L$, then $\mathfrak{A}$ has almost minimal spectrum (Theorem 2.1.4).
- If $\mathfrak{A}$ is a finite non-trivial algebra with almost minimal spectrum, then the variety generated by $\mathfrak{A}$ is congruence permutable (Theorem 2.2.4).
- If $\mathfrak{A}$ is a finite simple algebra in a permutable variety, then $\mathfrak{A}$ is either prime affine or its powers have only (trivial) factor congruences (Theorem 2.3.33).
- If $\mathfrak{A}$ is a finite simple non-trivial algebra with no proper subalgebras and no non-trivial automorphisms, and if $\mathfrak{A}$ generates a permutable and distributive variety, then $\mathfrak{A}$ is primal (Theorem 2.4.1).
2.0.10 Notation. Until the end of the chapter, as we will be proving Theorem 2.0.7, we will denote the algebra satisfying the hypotheses of the theorem by $\mathfrak{A}=(A, F)$. We will use the symbol $F$ also for the corresponding operations on powers of $\mathfrak{A}$. The congruence lattice of $\mathfrak{A}$ will play an important role and we will write $\operatorname{Con}(\mathfrak{A})$ for it. By $0 \in \operatorname{Con}(\mathfrak{A})$ we mean the diagonal $\{(a, a) \mid a \in A\}$ and by $1 \in \operatorname{Con}(\mathfrak{A})$ the trivial congruence $A^{2}$.


## $2.1 \mathfrak{A}$ has almost minimal spectrum

2.1.1 Definition. The $\operatorname{spectrum} \operatorname{Spec}(\mathcal{V})$ of a variety $\mathcal{V}$ is the set of all cardinalities of finite members of $\mathcal{V}$. For a finite algebra $\mathfrak{A}$ we define $\operatorname{Spec}(\mathfrak{A})=\operatorname{Spec}(\mathcal{V}(\mathfrak{A}))$, where $\mathcal{V}(\mathfrak{A})$ denotes the variety determined by $\mathfrak{A} . \mathfrak{A}$ is said to have minimal spectrum iff $\operatorname{Spec}(\mathfrak{A})=\left\{|\mathfrak{A}|^{n} \mid n \geq 0\right\}$.

The original goal of the author was to prove in this section that our algebra $\mathfrak{A}$ has minimal spectrum. This would have made it easy to find a title for this section. However, it did not work out and we will obtain that result later. The following definition will help us out for the moment.
2.1.2 Definition. We say that a finite algebra $\mathfrak{A}$ has almost minimal spectrum iff all subalgebras of finite powers of $\mathfrak{A}$ have cardinality a power of the cardinality of $\mathfrak{A}$.
2.1.3 Remark. Recall that every algebra in $\mathcal{V}(\mathfrak{A})$ is a homomorphic image of a subalgebra of a power of $\mathfrak{A}$. The notion of almost minimal spectrum is thus weaker than the one of minimal spectrum.

This section is devoted to the proof of the following theorem which is due to R. W. Quackenbush [9].
2.1.4 Theorem. Let $\mathfrak{A}$ be a finite non-trivial algebra having no subalgebra of a power of $\mathfrak{A}$ belonging to $R B L$. Then $\mathfrak{A}$ has almost minimal spectrum.

The proof will be by contradiction: Suppose $\mathfrak{A}$ does not have almost minimal spectrum; then there is an $m$ and a subalgebra $\mathfrak{B}$ of $\mathfrak{A}^{m}$ with $|B|$ not a power of $\kappa=|A|$. Choose $m$ minimal in the sense that for all $n<m$ every subalgebra of $\mathfrak{A}^{n}$ has cardinality a power of $\kappa$. As $\mathfrak{A}$ has no proper subalgebras (a proper subalgebra would be a unary central relation), $\mathfrak{B}$ must even be a subdirect product (that is, the projection of $\mathfrak{B}$ on any coordinate is onto); thus clearly, $m>1$.
2.1.5 Notation. For the rest of this section (that is, until Theorem 2.1.4 has been proven), we will extend Notation 2.0.10 and use the following conventions: $\mathfrak{A}$ will be assumed to satisfy all hypotheses of Theorem 2.1.4. The letter $\kappa$ will be reserved for the cardinality of $\mathfrak{A}$. For the universe of $\mathfrak{A}$ we write $A=\left\{\alpha_{1}, \ldots, \alpha_{\kappa}\right\} . \mathfrak{B}$ and $m$ as just defined will not change their meaning.

For $E=\left\{i_{1}, \ldots, i_{j}\right\} \subseteq\{1, \ldots, m\}$ where $i_{1}<\ldots<i_{j}$ define the projection

$$
\begin{array}{cccc}
\pi_{E}: & A^{m} & \rightarrow & A^{j} \\
\left(a_{1}, \ldots, a_{m}\right) & \mapsto & \left(a_{i_{1}}, \ldots, a_{i_{j}}\right) .
\end{array}
$$

Define further for $1 \leq i \leq m$ the projection $\rho_{i}=\pi_{E(i)}$ where $E(i)=\{1, \ldots, i-1, i+1, \ldots, m\}$.
2.1.6 Lemma. $\left|\rho_{i}(B)\right|=\kappa^{m-1}$ for $1 \leq i \leq m$.

Proof. We will proof by induction on $n$ that any projection $\pi$ mapping from $B$ to $A^{n}$, where $n<m$, is onto. Since $\mathfrak{A}$ has no proper subalgebras, all the projections $\pi_{i}^{m}$ from $B$ to $A$ are onto and hence our assertion is true for $n=1$. Now assume that for all $i<n<m$, if $|E|=i$ then $\left|\pi_{E}(B)\right|=|A|^{i}$ and let $E$ be an $n$-element subset of $\{1, \ldots, m\}$. As $\pi_{E}(B)$ is a subalgebra of $A^{n}$ and $n<m$, there exists an $i$ such that $\left|\pi_{E}(B)\right|=|A|^{i}$. Trivially, $\left|\pi_{E}(B)\right| \leq|A|^{n}$ and by induction hypothesis, $\left|\pi_{E}(B)\right| \geq|A|^{n-1}$ so that either $\left|\pi_{E}(B)\right|=|A|^{n-1}$ or $\left|\pi_{E}(B)\right|=|A|^{n}$. Consider in the first case $E^{\prime}=E \backslash\{i\}$ for an arbitrary $i \in E$. By induction hypothesis we know that also $\left|\pi_{E^{\prime}}(B)\right|=|A|^{n-1}$. From this follows that for $b, b^{\prime} \in B$, if $b_{j}=b_{j}^{\prime}$ for $j \in E^{\prime}$ then $b_{i}=b_{i}^{\prime}$. Hence for any $b, b^{\prime} \in B, \rho_{i}(b)=\rho_{i}\left(b^{\prime}\right)$ implies that $b=b^{\prime}$ and $\rho_{i}$ is one-one. But this means that $\rho_{i}$ embeds $\mathfrak{B}$ as a subalgebra of $\mathfrak{A}^{m-1}$, contradicting our assumption that $m$ is minimal with respect to having a subalgebra of cardinality not a power of $A$. Thus $\left|\pi_{E}(B)\right|$ must be equal to $|A|^{n}$ and the induction is complete.
2.1.7 Corollary. $|A|^{m-1}<|B|<|A|^{m}$.

Let $P$ be a partition $\{1, \ldots, m\}$, and let $\sim_{P}$ be the equivalence relation induced by $P$. Define a subset $B_{P}$ of $B$ by $B_{P}=\left\{b \in B \mid i \sim_{P} j \rightarrow b_{i}=b_{j}\right\}$. Then clearly, $B_{P}$ is a subuniverse of $B$. To denote a partition, we will only list its non-trivial classes; $(i, j)$ denotes the partition with only one non-trivial class, $\{i, j\}$.
2.1.8 Lemma. Let $m \geq 4$. If for some $1 \leq i, j \leq m, i \neq j$ we have $\left|B_{(i, j)}\right|=|A|^{m-1}$, then the same holds for all $1 \leq i, j \leq m, i \neq j$.

Proof. Let $\left|B_{(i, j)}\right|=|A|^{m-1}$; it suffices to show that for $k \neq i, j$ we have $\left|B_{(i, k)}\right|=|A|^{m-1}$. Our assumption $\left|B_{(i, j)}\right|=|A|^{m-1}$ obviously implies $\left|B_{(i, j, k)}\right|=|A|^{m-2}$. Since $B_{(i, k)}$ can be embedded into $\mathfrak{A}^{m-1}$ by leaving away the $k$-th coordinate, $\left|B_{(i, k)}\right|$ must be a power of $|A|$. Trivially, $\left|B_{(i, k)}\right| \leq|A|^{m-1}$ and since $\left|B_{(i, j, k)}\right| \leq\left|B_{(i, k)}\right|$ we have $\left|B_{(i, k)}\right| \geq|A|^{m-2}$. Suppose now that $\left|B_{(i, k)}\right|=|A|^{m-2}$; then $\left|B_{(i, j, k)}\right|=\left|B_{(i, k)}\right|$ and so $b_{i}=b_{k}$ implies $b_{i}=b_{j}$ for all $b \in B$. But this contradicts that by the proof of Lemma 2.1.6, $\left|\pi_{\{i, j, k\}}(B)\right|=|A|^{3}$. Therefore, $\left|B_{(i, k)}\right| \neq|A|^{m-2}$ and so $\left|B_{(i, k)}\right|=|A|^{m-1}$.

Define a subset $B^{\prime}$ of $A^{m}$ by

$$
B^{\prime}=\left\{\left(a_{2}, a_{3}, \ldots, a_{m}, a_{m}^{\prime}\right) \mid \exists a_{1} \in A\left(\left(a_{1}, \ldots, a_{m}\right) \in B \wedge\left(a_{1}, \ldots, a_{m-1}, a_{m}^{\prime}\right) \in B\right)\right\}
$$

Then $B^{\prime}$ is a subuniverse of $\mathfrak{A}^{m}$ and the following holds:
2.1.9 Lemma. Let $m \geq 4$. If $\left|B_{(2,3)}\right|=|A|^{m-2}$, then

- $|A|^{m-1}<\left|B^{\prime}\right|<|A|^{m}$
- $\left|B_{(1,2)}^{\prime}\right|=|A|^{m-2}$
- $\left|B_{(m-1, m)}^{\prime}\right|=|A|^{m-1}$

Proof. First note that $\left|B_{(m-1, m)}^{\prime}\right|=\left|\rho_{1}(B)\right|=|A|^{m-1}$, the latter equality provided by Lemma 2.1.6. Furthermore, $\left|B^{\prime}\right| \geq\left|B_{(m-1, m)}^{\prime}\right|=|A|^{m-1}$. Since we know that $\left|\rho_{m}(B)\right|=|A|^{m-1}$ but $|B|>|A|^{m-1}$, there exist $a_{1}, \ldots, a_{m-1}, a_{m}, a_{m}^{\prime} \in A$ such that $a_{m} \neq a_{m}^{\prime}$ and both $\left(a_{1}, \ldots, a_{m-1}, a_{m}\right) \in B$ and $\left(a_{1}, \ldots, a_{m-1}, a_{m}^{\prime}\right) \in B$. Hence, $\left(a_{2}, \ldots, a_{m-1}, a_{m}, a_{m}^{\prime}\right) \in B^{\prime}$ and so $\left|B^{\prime}\right|>|A|^{m-1}$. Given $a_{1}, \ldots, a_{m-1} \in A,\left|\rho_{m}(B)\right|=|A|^{m-1}$ implies there exists an $a_{m} \in A$ such that $\left(a_{1}, \ldots, a_{m}\right) \in B$. Since we assume $\left|B_{(2,3)}\right|=|A|^{m-2}$, from $a_{2}=a_{3}$ it follows that such an $a_{m}$ is unique so that if we choose any $a_{m}^{\prime} \neq a_{m},\left(a_{3}, a_{3}, a_{4}, \ldots, a_{m-1}, a_{m}, a_{m}^{\prime}\right) \notin B^{\prime}$. Thus, $\left|B^{\prime}\right|<|A|^{m}$, and $\left|B_{(1,2)}^{\prime}\right|=\left|B_{(2,3)}\right|=|A|^{m-2}$.
2.1.10 Lemma. Let $m \geq 4$. Then for $1 \leq i<j \leq m,\left|B_{(i, j)}\right|=|A|^{m-1}$.

Proof. If $\left|B_{(2,3)}\right|=|A|^{m-1}$ then the assertion follows from Lemma 2.1.8. If not, then $\left|B_{(2,3)}\right|=$ $|A|^{m-2}$, and so by the last lemma $\left|B_{(m-1, m)}^{\prime}\right|=|A|^{m-1}$. Thus, $B^{\prime}$ satisfies the hypotheses on $B$ in Lemmas 2.1.6 and 2.1.8 and application of Lemma 2.1.8 yields $\left|B_{(1,2)}^{\prime}\right|=|A|^{m-1}$ contradicting $\left|B_{(1,2)}^{\prime}\right|=|A|^{m-2}$ which we established in the previous lemma. Hence, $\left|B_{(2,3)}\right|=|A|^{m-2}$ is impossible and the lemma follows.

We can summarize what we have established so far:
2.1.11 Theorem. Let $\mathfrak{B}$ be a subalgebra of $\mathfrak{A}^{m}$ with $|B|$ not a power of $\kappa$. If $m \geq 4$, then $B$ is totally reflexive.

Denote by $\mathfrak{B}^{*}=\left(B^{*}, F\right)$ the subalgebra of $\mathfrak{A}^{m}$ generated by $\iota_{m}^{A}$. By considering this algebra we will show that for $m \geq 4$ we can assume without loss of generality that $\mathfrak{B}$ is totally reflexive and totally symmetric:
2.1.12 Theorem. Let $m \geq 4$. Then $B^{*}$ is totally reflexive and totally symmetric and $|A|^{m-1}<$ $\left|B^{*}\right|<|A|^{m}$.

Proof. $\iota_{m}^{A}$ is both totally reflexive and totally symmetric and it is easy to see that $B^{*}$ inherits those properties. By Theorem 2.1.11 we have $\iota_{m}^{A} \subseteq B$ and so $B^{*} \subseteq B$; hence, $\left|B^{*}\right| \leq|B|<|A|^{m}$. Moreover, $\left|B_{(1,2)}\right|=|A|^{m-1}$ by Lemma 2.1.10, and $B_{(1,2)}$ is obviously a proper subset of $\iota_{m}^{A}$. Thus, $|A|^{m-1}<\left|\iota_{m}^{A}\right|<\left|B^{*}\right|$ and the theorem follows.

We have shown that in the case $m \geq 4$, we can assume $B$ to be totally reflexive and totally symmetric by replacing $B$ with $B^{*}$ if necessary. Our next step will be to prove the totally reflexive and totally symmetric possibility absurd; as a result, $m \geq 4$ cannot occur.

## The totally reflexive and totally symmetric case

First note that in this case $m \leq \kappa$ since otherwise every element of $A^{m}$ would have two equal components and so the total reflexivity of $B$ would imply $B=A^{m}$. Now choose $h \leq \kappa$ to be maximal with respect to $\mathfrak{A}^{h}$ containing a proper totally reflexive and totally symmetric subalgebra; let $\mathfrak{C}=(C, F)$ be a maximal subalgebra of $\mathfrak{A}^{h}$ of that kind.

For $h \leq n \leq \kappa$ define sets $C_{n} \subseteq A^{n}$ to contain all $\left(a_{1}, \ldots, a_{n}\right)$ for which there exists an $a \in A$ such that for each $(h-1)$-element subset $\left\{i_{1}, \ldots, i_{h-1}\right\}$ of $\{1, \ldots, n\},\left(a_{i_{1}}, \ldots, a_{i_{h-1}}, a\right) \in C$. Then $\left(C_{n}, F\right)$ is a subalgebra of $\mathfrak{A}^{n}$ and is totally symmetric as $C$ is.
2.1.13 Lemma. Either $C=C_{h}$ or $C_{h}=A^{h}$.

Proof. Let $\left(a_{1}, \ldots, a_{h}\right) \in C$. Set $a=a_{1}$; then by the total reflexivity and total symmetry of $C$ we have that for $1 \leq i \leq h,\left(a_{1}, \ldots, a_{i-1}, a_{i+1}, \ldots, a_{h}, a_{1}\right) \in C$ so that $\left(a_{1}, \ldots, a_{h}\right) \in C_{h}$. Hence, $C \subseteq C_{h}$ and so by the maximality of $C, C=C_{h}$ or $C_{h}=A^{h}$.
2.1.14 Lemma. If $C_{h}=A^{h}$, then $C_{\kappa}=A^{\kappa}$.

Proof. All $C_{n}$ are totally symmetric, $h \leq n \leq \kappa$. Thus, by the maximality of $h$, if $C_{n}$ is also totally reflexive then $C_{n}=A^{n}$. But clearly the definition of $C_{n}$ implies that if $C_{n}=A^{n}$, then $C_{n+1}$ is totally reflexive so that by induction we get $C_{\kappa}=A^{\kappa}$.

The following lemma states that the case $C_{h}=A^{h}$ is impossible.
2.1.15 Lemma. If $C \neq C_{h}$, then $C$ is central.

Proof. By Lemmas 2.1.13 and 2.1.14 our hypothesis implies that $C_{\kappa}=A^{\kappa}$. Hence, $\left(\alpha_{1}, \ldots, \alpha_{\kappa}\right) \in C_{\kappa}$. Therefore, by definition of $C_{\kappa}$, there is an $a \in A$ such that for every $h-1$-element subset $\left\{\alpha_{i_{1}}, \ldots, \alpha_{i_{h-1}}\right\}$ of $\left\{\alpha_{1}, \ldots, \alpha_{\kappa}\right\}=A,\left(\alpha_{i_{1}}, \ldots, \alpha_{i_{h-1}}, a\right) \in C$. Together with the fact that $C$ is totally reflexive and totally symmetric this implies that $a$ is in the center of $C$. But since $C \neq A^{h}$ this means that $C$ is central.

We may therefore assume that $C=C_{h}$. Such a $C$ is called homogeneous. Note that in the case $h=2$, if $\left(a_{1}, a_{3}\right) \in C$ and $\left(a_{2}, a_{3}\right) \in C$, we have that also $\left(a_{1}, a_{2}\right) \in C$ (set $a=a_{3}$ in the definition of $C_{h}$ ). As $C$ is symmetric this means that $C$ is transitive and so, together with its reflexivity, we get that $C$ is a non-trivial equivalence relation and thus in Rosenberg's list. Therefore, we may assume $h \geq 3$.

In the following, we will make use of the homogeneity of $C$. Notice therefore that by the definition of $C_{h}$, to prove that a tuple $\left(a_{1}, \ldots, a_{h}\right)$ it is an element of $C$ it suffices to find an arbitrary $b \in A$ such that for all $1 \leq i \leq h$, if we replace $a_{i}$ by $b$, then the resulting tuple is in $C$. Such an element will be referred to as a replacement element. The condition is not only sufficient but also necessary for membership of $C$.

If $C$ contains all tuples $\left(a_{1}, \ldots, a_{h}\right) \in A^{h}$ for which there exists $\left(v_{1}, \ldots, v_{h}\right) \in C$ such that $\left(a_{1}, \ldots, a_{i-1}, a_{i+1}, \ldots, a_{h}, v_{j}\right) \in C$ for all $1 \leq i, j \leq h$ with $i \neq j$, then it is called strongly homogeneous. The tuple $\left(v_{1}, \ldots, v_{h}\right)$ will be referred to as the replacement tuple. Notice that strongly homogeneous immediately implies homogeneous if one considers the replacement tuple containing a replacement element at every coordinate. We will see that the maximality of $C$ implies that it is strongly homogeneous; from that we will derive that $C$ is $h$-regularly generated so that it belongs to Rosenberg's list.

For $h \leq r \leq \kappa$, define $C_{r} \subseteq A^{r}$ by

$$
C_{r}=\left\{\left(a_{1}, \ldots, a_{r}\right) \mid \forall E \subseteq\{1, \ldots, r\}\left(|E|=h \rightarrow \pi_{E}\left(a_{1}, \ldots, a_{r}\right) \in C\right)\right\}
$$

Clearly for all $r,\left(C_{r}, F\right)$ is a subalgebra of $\mathfrak{A}^{r}, C_{r}$ is totally symmetric since $C$ is, and $C_{h}=C$. For $h \leq r \leq \kappa$, define $D_{r} \subseteq A^{r}$ by

$$
\begin{aligned}
D_{r}= & \left\{\left(a_{1}, \ldots, a_{r}\right) \mid \exists\left(b_{1}, \ldots, b_{r}\right) \in C_{r}\right. \\
& \left.\forall 1 \leq j \leq r \forall\left\{i_{1}, \ldots, i_{h-2}\right\} \subseteq\{1, \ldots, r\} \quad\left(a_{i_{1}}, \ldots, a_{i_{h-2}}, a_{j}, b_{j}\right) \in C\right\}
\end{aligned}
$$

Then for all $r,\left(D_{r}, F\right)$ is a subalgebra of $\mathfrak{A}^{r}$; furthermore, $D_{r}$ is totally symmetric by its symmetric definition and the total symmetry of $C_{r}$.
2.1.16 Lemma. If $C=D_{h}$ then $C$ is strongly homogeneous.

Proof. Suppose $C=D_{h}$ and let $\left(a_{1}, \ldots, a_{h}\right) \in A^{h}$ and $\left(v_{1}, \ldots, v_{h}\right) \in C$ be given such that $\left(a_{1}, \ldots, a_{i-1}, a_{i+1}, \ldots, a_{h}, v_{j}\right) \in C$ for all $1 \leq i, j \leq h$ with $i \neq j$. Then for every $1 \leq j \leq h$ and every $\left\{i_{1}, \ldots, i_{h-2}\right\} \subseteq\{1, \ldots, h\}$ we have that $\left(a_{i_{1}}, \ldots, a_{i_{h-2}}, a_{j}, v_{j}\right) \in C$. This is clear if $\left|\left\{i_{1}, \ldots, i_{h-2}, j\right\}\right|<h-1$ from the total reflexivity of $C$ and if not, then there is an $i \neq j$ such that $\left(a_{i_{1}}, \ldots, a_{i_{h-2}}, a_{j}, v_{j}\right) \in C$ is by the total symmetry of $C$ equivalent to $\left(a_{1}, \ldots, a_{i-1}, a_{i+1}, \ldots, a_{h}, v_{j}\right) \in C$. Since the latter statement is true, setting $\left(b_{1}, \ldots, b_{h}\right)=\left(v_{1}, \ldots, v_{h}\right)$ in the definition of $D_{h}$ shows $\left(a_{1}, \ldots, a_{h}\right) \in D_{h}=C$. Hence, $C$ is indeed strongly homogeneous.
2.1.17 Lemma. Either $C=D_{h}$ or $D_{h}=A^{h}$.

Proof. Let $\left(a_{1}, \ldots, a_{h}\right) \in C$ and set in the definition of $D_{h}\left(b_{1}, \ldots, b_{h}\right)$ equal to $\left(a_{1}, \ldots, a_{h}\right)$. Then $\left(b_{1}, \ldots, b_{h}\right) \in C_{h}=C$, and so $\left(a_{1}, \ldots, a_{h}\right) \in D_{h}$ as $C$ is totally reflexive. Therefore, $C \subseteq D_{h}$ and consequently the maximality of $h$ implies $C=D_{h}$ or $D_{h}=A^{h}$.
2.1.18 Lemma. If $D_{h}=A^{h}$, then $D_{\kappa}=A^{\kappa}$.

Proof. The proof will be by induction. Suppose that $D_{n}=A^{n}$; choose an arbitrary $\left(a_{1}, \ldots, a_{n}\right) \in A^{n}$ and say it is in $D_{n}$ via $\left(b_{1}, \ldots, b_{n}\right) \in C_{n}$. But if $\left(b_{1}, \ldots, b_{n}\right) \in C_{n}$, then obviously $\left(b_{1}, b_{1}, \ldots, b_{n}\right) \in C_{n+1}$ and it is easily seen that $\left(a_{1}, a_{1}, \ldots, a_{n}\right) \in D_{n+1}$ via $\left(b_{1}, b_{1}, \ldots, b_{n}\right)$. Since we could have chosen any other coordinate instead of the first in that argument, we conclude that $D_{n+1}$ is totally reflexive. But $D_{n+1}$ is also totally symmetric, and so the maximality of $C$ implies $D_{n+1}=A^{n+1}$.
2.1.19 Theorem. $C$ is strongly homogeneous.

Proof. We will show that $D_{\kappa} \neq A^{\kappa}$. Then by Lemma $2.1 .18, D_{h} \neq A^{h}$, and so by Lemma 2.1.17, $C=D_{h}$ which we know implies that $C$ is strongly homogeneous. Suppose towards contradiction that $D_{\kappa}=A^{\kappa}$ : then the vector $\left(a_{1}, \ldots, a_{\kappa}\right) \in A^{\kappa}$ that lists all elements of $A$ is an element of $D_{\kappa}$, say via $\left(b_{1}, \ldots, b_{\kappa}\right) \in C_{\kappa}$. By definition of $D_{\kappa}$, for every $1 \leq j \leq \kappa$ and every $\left\{i_{1}, \ldots, i_{h-2}\right\} \subseteq\{1, \ldots, \kappa\}$ we have $\left(a_{i_{1}}, \ldots, a_{i_{h-2}}, a_{j}, b_{j}\right) \in C$. Since $\left(a_{1}, \ldots, a_{\kappa}\right)$ lists $A$, this means that all $h$-tuples containing $a_{j}$ and $b_{j}$ for some $j$ are in $C$. We will prove by induction that for $0 \leq n \leq h,\left(a_{1}, \ldots, a_{n}, b_{n+1}, \ldots, b_{h}\right) \in C$. Since $\left(b_{1}, \ldots, b_{\kappa}\right) \in C_{\kappa},\left(b_{1}, \ldots, b_{h}\right) \in C_{h}=C$ and so in the case $n=0$ our assertion is true. Suppose it is true for $n<h$ and consider $\left(a_{1}, \ldots, a_{n+1}, b_{n+2}, \ldots, b_{h}\right)$. If we replace any element other that $a_{n+1}$ by $b_{n+1}$, the resulting $h$-tuple contains both $a_{n+1}$ and $b_{n+1}$ so that it is in $C$ by the preceding discussion. On the other hand, replacing $a_{n+1}$ by $b_{n+1}$ gives us an element of $C$ by induction hypothesis. Hence, the homogeneity of $C$ implies that $\left(a_{1}, \ldots, a_{n+1}, b_{n+2}, \ldots, b_{h}\right) \in C$ and the induction is complete. Now setting $n=h$ yields $\left(a_{1}, \ldots, a_{h}\right) \in C$. But the vector $\left(a_{1}, \ldots, a_{\kappa}\right)$ was arbitrarily chosen; hence, $C=A^{h}$, contradicting our assumption on $C$.

## The case of the $h$-regularly generated relations

We will show that strongly homogeneous $C$ is $h$-regularly generated. That is, we will find a surjection $\varphi: A \rightarrow h^{m}$ such that $C=\varphi^{-1}\left(\omega_{m}\right)$ as defined in the first chapter. Our first step is to find the equivalence relation induced by $\varphi$. Define $E \subseteq A^{2}$ by

$$
E=\left\{(a, b) \mid \forall\left(a_{1}, \ldots, a_{h-2}\right) \in A^{h-2}\left(a_{1}, \ldots, a_{h-2}, a, b\right) \in C\right\} ;
$$

then $E$ is an equivalence relation on $A$. Reflexivity and symmetry of $E$ immediately follow from the corresponding properties of $C$. To see $E$ is transitive, let $(a, b),(b, c) \in E$ be given. Then since $C$ is homogeneous, using $b$ as a replacement element yields that for all $\left(a_{1},, \ldots, a_{h-2}\right) \in$ $A^{h-2},\left(a_{1}, \ldots, a_{h-2}, a, c\right) \in C$ and thus $(a, c) \in E$. Suppose now that $E$ has $q$ equivalence classes and assume without loss of generality that $A^{\prime}=\left\{\alpha_{1}, \ldots, \alpha_{q}\right\}$ contains one element from each equivalence class. Let $\gamma: A \rightarrow A^{\prime}$ be the function that maps each $a \in A$ to the element in $A^{\prime}$ that represents the equivalence class of $a$; that is, $(a, \gamma(a)) \in E$ for all $a \in A$. Define $C^{*} \subseteq A^{h}$ by

$$
C^{*}=\left\{\left(a_{1}, \ldots, a_{h}\right) \mid\left(\gamma\left(a_{1}\right), \ldots, \gamma\left(a_{h}\right)\right) \in C\right\} .
$$

2.1.20 Lemma. Let $(a, b) \in E$ and $\left(a_{1}, \ldots, a_{h-1}\right) \in A^{h-1}$. Then $\left(a_{1}, \ldots, a_{h-1}, a\right) \in C$ iff $\left(a_{1}, \ldots, a_{h-1}, b\right) \in C$.

Proof. Let $\left(a_{1}, \ldots, a_{h-1}, a\right) \in C$. By definition of $E$ and the total symmetry of $C$, $\left(a_{1}, \ldots, a_{i-1}, a, a_{i+1}, \ldots, a_{h-1}, b\right) \in C$ for all $1 \leq i \leq h-1$. Thus, if we use $a$ as a replacement element, the homogeneity of $C$ implies $\left(a_{1}, \ldots, a_{h-1}, b\right) \in C$.
2.1.21 Theorem. $C^{*}=C$. That is, membership of $C$ is completely determined by the equivalence classes of $E$.

Proof. First, let $\left(a_{1}, \ldots, a_{h}\right) \in C$. Then by the previous lemma, $\left(\gamma\left(a_{1}\right), a_{2}, \ldots, a_{h}\right) \in C$. Hence by induction, $\left(\gamma\left(a_{1}\right), \gamma\left(a_{2}\right), \ldots, \gamma\left(a_{h}\right)\right) \in C$ so that $\left(a_{1}, \ldots, a_{h}\right) \in C^{*}$. Conversely, if $\left(a_{1}, \ldots, a_{h}\right) \in C^{*}$, then $\left(\gamma\left(a_{1}\right), \gamma\left(a_{2}\right), \ldots, \gamma\left(a_{h}\right)\right) \in C$, and applying the same induction backwards yields $\left(a_{1}, \ldots, a_{h}\right) \in C$.

Call $C$ universal if there exists a function $f: h^{h^{\kappa}} \rightarrow A$ such that for $1 \leq j \leq \kappa, f\left(\pi_{j}^{\kappa}\right)=\alpha_{j}$ and such that for all $\left(b_{1}, \ldots, b_{h}\right) \in \omega_{h^{\kappa}},\left(f\left(b_{1}\right), \ldots, f\left(b_{h}\right)\right) \in C$. Our next goal is to prove that $C$ is universal. For $h \leq i \leq \kappa$ define $\bar{C}_{i} \subseteq A^{i}$ by $\left(a_{1}, \ldots, a_{i}\right) \in \bar{C}_{i}$ iff there is an $f: h^{h^{i}} \rightarrow A$ such that for $1 \leq j \leq i, f\left(\pi_{j}^{i}\right)=a_{j}$ and such that for all $\left(b_{1}, \ldots, b_{h}\right) \in \omega_{h^{i}},\left(f\left(b_{1}\right), \ldots, f\left(b_{h}\right)\right) \in C$. We will prove that $\bar{C}_{\kappa}=A^{\kappa}$ to show that $C$ is universal.
2.1.22 Lemma. $\bar{C}_{h}=A^{h}$.

Proof. Let $\left(a_{1}, \ldots, a_{h}\right) \in A^{h}$ be given and let $z=(0, \ldots, h-1) \in h^{h}$. Define $\omega: h \rightarrow A$ by $\omega(j)=a_{j+1}$ for $0 \leq j \leq h-1$ and define $f_{\omega}: h^{h^{h}} \rightarrow A$ by $f_{\omega}(b)=\omega(b(z))$ for all $b \in h^{h^{h}}$. Then $f_{\omega}\left(\pi_{j}^{h}\right)=\omega\left(\pi_{j}^{h}(z)\right)=\omega(j-1)=a_{j}$ for $1 \leq j \leq h$. Moreover, if $\left(b_{1}, \ldots, b_{h}\right) \in \omega_{h^{h}}$,
then $\left(f_{\omega}\left(b_{1}\right), \ldots, f_{\omega}\left(b_{h}\right)\right)=\left(\omega\left(b_{1}(z)\right), \ldots, \omega\left(b_{h}(z)\right)\right) \in \iota_{h}^{A} \subseteq C$ since $\left(b_{1}(z), \ldots, b_{h}(z)\right) \in \iota_{h}^{h}$ by definition of $\omega_{h^{h}}$. Thus, $\left(a_{1}, \ldots, a_{h}\right) \in C_{h}$.
2.1.23 Lemma. Let $h \leq i \leq \kappa$. Then $\left(\bar{C}_{i}, F\right)$ is a totally symmetric subalgebra of $\mathfrak{A}^{i}$.

Proof. Let $g$ be an $n$-ary operation of $\mathfrak{A}$ and for $1 \leq j \leq n$ let $a_{j} \in \bar{C}_{i}$ via the function $f_{j}: h^{h^{i}} \rightarrow A$. Set $f=g\left(f_{1}, \ldots, f_{n}\right)$ and write $g\left(a_{1}, \ldots, a_{n}\right)=\left(d_{1}, \ldots, d_{i}\right) \in A^{i}$. We will show that $d=\left(d_{1}, \ldots, d_{i}\right) \in \bar{C}_{i}$ via $f$. First note that if $a_{j}=\left(a_{j 1}, \ldots, a_{j i}\right)$, then $f\left(\pi_{j}^{i}\right)=g\left(f_{1}\left(\pi_{j}^{i}\right), \ldots, f_{n}\left(\pi_{j}^{i}\right)\right)=g\left(a_{1 j}, \ldots, a_{n j}\right)=d_{j}$. Moreover, if $\left(b_{1}, \ldots, b_{h}\right) \in \omega_{h^{i}}$, then $\left(f\left(b_{1}\right), \ldots, f\left(b_{h}\right)\right)=\left(g\left(f_{1}\left(b_{1}\right), \ldots, f_{n}\left(b_{1}\right)\right), \ldots, g\left(f_{1}\left(b_{h}\right), \ldots, f_{n}\left(b_{h}\right)\right)\right) \in C$ since already $\left(f_{j}\left(b_{1}\right), \ldots, f_{j}\left(b_{h}\right)\right) \in C$ for $1 \leq j \leq n$ and since $\mathfrak{C}$ is a subalgebra of $\mathfrak{A}^{h}$. Hence, $d$ is indeed an element of $\bar{C}_{i}$ via $f$ and thus $\left(\bar{C}_{i}, F\right)$ is a subalgebra of $\mathfrak{A}^{i}$. To show that $\bar{C}_{i}$ is totally symmetric, let $\sigma$ be any permutation of $\{1, \ldots, i\}$ and let $\left(a_{1}, \ldots, a_{i}\right) \in \bar{C}_{i}$ via $f$. Then $\left(a_{\sigma(1)}, \ldots, a_{\sigma(i)}\right) \in \bar{C}_{i}$ via $f_{\sigma}$ if we set $f_{\sigma}(b)=f(\tilde{b})$ where $\tilde{b}\left(x_{1}, \ldots, x_{i}\right)=b\left(x_{\sigma(1)}, \ldots, x_{\sigma(i)}\right)$. For $f_{\sigma}\left(\pi_{j}^{i}\right)=f\left(\pi_{\sigma(j)}^{i}\right)=a_{\sigma(j)}$ and if $\left(b_{1}, \ldots, b_{h}\right) \in \omega_{h^{i}}$, then also ( $\left.\tilde{b_{1}}, \ldots, \tilde{b_{h}}\right) \in \omega_{h^{i}}$ so that $\left(f_{\sigma}\left(b_{1}\right), \ldots, f_{\sigma}\left(b_{h}\right)\right)=\left(f\left(\tilde{b_{1}}\right), \ldots, f\left(\tilde{b_{h}}\right)\right) \in C$.
2.1.24 Lemma. $C$ is universal.

Proof. We will prove by induction that for $h \leq n \leq \kappa, \bar{C}_{n}=A^{n}$. By Lemma 2.1.22, $\bar{C}_{h}=A^{h}$. Now assume $\bar{C}_{n}=A^{n}$; we will show that this implies that $\bar{C}_{n+1}$ is totally reflexive. Let $\left(a_{1}, \ldots, a_{n}\right) \in A^{n}=\bar{C}_{n}$ via $f$ and define $\tilde{f}: h^{h^{n+1}} \rightarrow A$ by $\tilde{f}\left(b\left(x_{1}, \ldots, x_{n+1}\right)\right)=$ $f\left(b\left(x_{2}, x_{2}, \ldots, x_{n+1}\right)\right)$. Then it is easy to see that $\left(a_{1}, a_{1}, a_{2}, \ldots, a_{n}\right) \in \bar{C}_{n+1}$ so that since $\left(a_{1}, \ldots, a_{n}\right)$ was an arbitrary tuple in $A^{n}$ we have that $\bar{C}_{n+1}$ is totally reflexive and must therefore equal $A^{n+1}$. Now in particular $\bar{C}_{\kappa}=A^{\kappa}$ and hence, $\left(\alpha_{1}, \ldots, \alpha_{\kappa}\right) \in \bar{C}_{\kappa}$ which means exactly that $C$ is universal.

In the light of Theorem 2.1.21, it is natural to consider $D=C \cap\left(A^{\prime}\right)^{h} . D$ is totally reflexive, totally symmetric and strongly homogeneous, the latter since Lemma 2.1.20 implies that in the definition of strong homogeneity, we can replace $\left(v_{1}, \ldots, v_{h}\right)$ with $\left(\gamma\left(v_{1}\right), \ldots, \gamma\left(v_{h}\right)\right)$. Fix $f^{\prime}: h^{h^{\kappa}} \rightarrow A$ making $C$ universal and set $f=\gamma \circ f^{\prime}: h^{h^{\kappa}} \rightarrow A^{\prime} . f$ makes $D$ kind of universal in the sense that for all $\left(b_{1}, \ldots, b_{h}\right) \in \omega_{h^{\kappa}},\left(f\left(b_{1}\right), \ldots, f\left(b_{h}\right)\right) \in D$; this is a consequence of Lemma 2.1.21. However, $D$ is not a subuniverse of $\mathfrak{A}^{h}$ as $A^{\prime}$ is not closed under those operations.

We will prove that there is an $\lambda$ such that $q=\left|A^{\prime}\right|=h^{\lambda}$. Let $s, t \in h^{\kappa}, j \in h, g: h^{\kappa} \rightarrow h$. Define $g_{j}^{t}: h^{\kappa} \rightarrow h$ by

$$
g_{j}^{t}(s)= \begin{cases}g(s) & , s \neq t \\ j & , \text { otherwise }\end{cases}
$$

Define $B_{j}^{t}=\left\{b \in h^{h^{\kappa}} \mid b(t)=j\right\}$.
2.1.25 Lemma. Let $\left(a_{1}, \ldots, a_{h}\right) \in \iota_{h}^{h}, t \in h^{\kappa}, g: h^{\kappa} \rightarrow h$, and $b_{i} \in B_{a_{i}}^{t}$ for $1 \leq i \leq h-2$. Then $\left(f\left(b_{1}\right), \ldots, f\left(b_{h-2}\right), f\left(g_{a_{h-1}}^{t}\right), f\left(g_{a_{h}}^{t}\right)\right) \in D$.

Proof. Since $D$ is kind of universal via $f$, it is enough to show that $\left(b_{1}, \ldots, b_{h-2}, g_{a_{h-1}}^{t}, g_{a_{h}}^{t}\right) \in$ $\omega_{h^{\kappa}}$. Let $t^{\prime} \in h^{\kappa}$ be given; we evaluate the tuple above at $t^{\prime}$. If $t^{\prime}=t$, then we get $\left(a_{1}, \ldots, a_{h-2}, a_{h-1}, a_{h}\right) \in \iota_{h}^{h}$; if $t^{\prime} \neq t$, then we get $\left(b_{1}\left(t^{\prime}\right), \ldots, b_{h-2}\left(t^{\prime}\right), g\left(t^{\prime}\right), g\left(t^{\prime}\right)\right) \in \iota_{h}^{h}$. Therefore $\left(b_{1}, \ldots, b_{h-2}, g_{a_{h-1}}^{t}, g_{a_{h}}^{t}\right) \in \omega_{h^{\kappa}}$ by definition of $\omega_{h^{\kappa}}$.
2.1.26 Lemma. Let $g: h^{\kappa} \rightarrow h, t \in h^{\kappa}, b_{p}: h^{\kappa} \rightarrow h$ for $1 \leq p \leq h$, and $\left(f\left(g_{0}^{t}\right), \ldots, f\left(g_{h-1}^{t}\right)\right) \in$ D. Consider $\left(f\left(b_{1}\right), \ldots, f\left(b_{h}\right)\right)$; then for all $1 \leq r<s \leq h$ the following holds: If we replace the $r$-th component, $f\left(b_{r}\right)$, by $f\left(g_{r-1}^{t}\right)$ and the s-th component, $f\left(b_{s}\right)$, by $f\left(g_{s-1}^{t}\right)$, then the resulting tuple is an element of $D$.

Proof. Since the $b_{p}$ are arbitrarily given and since $D$ is totally symmetric, it suffices to consider $r=h-1$ and $s=h$. If $\left(b_{1}(t), \ldots, b_{h-2}(t), h-2, h-1\right) \in \iota_{h}^{h}$, then by setting $\left(a_{1}, \ldots, a_{h}\right)=$ $\left(b_{1}(t), \ldots, b_{h-2}(t), h-2, h-1\right)$ the result follows from the previous lemma. Otherwise, assume without loss of generality that for $1 \leq p \leq h-2, b_{p}(t)=p-1$. We will proof by induction that for $0 \leq n \leq h-2,\left(f\left(b_{1}\right), \ldots, f\left(b_{n}\right), f\left(g_{n}^{t}\right), \ldots, f\left(g_{h-1}^{t}\right)\right) \in D$. For $n=0$, this is an assumption of the lemma; suppose it holds for $n<h-2$. Consider $\left(b_{1}, \ldots, b_{n+1}, g_{n+1}^{t}, \ldots, g_{h-1}^{t}\right)$. Replace any component other than $b_{n+1}$ by $g_{n}^{t}$. We have that $g_{n}^{t}(t)=n=b_{n+1}(t)$, whereas for all $t^{\prime} \neq t, g_{n}^{t}\left(t^{\prime}\right)=g_{h-2}^{t}\left(t^{\prime}\right)=g_{h-1}^{t}\left(t^{\prime}\right)=g\left(t^{\prime}\right)$. Thus, each of the tuples that result from our replacement belongs to $\omega_{h^{\kappa}}$ so that since $D$ is kind of universal, applying $f$ to each coordinate of such replacement tuples results in a member of $D$. By induction hypothesis, replacing $b_{n+1}$ by $g_{n}^{t}$ and application of $f$ yields a member of $D$ too, so that using the homogeneity of $D$ with $g_{n}^{t}$ as the replacement element concludes the proof.
2.1.27 Lemma. Suppose all assumptions of Lemma 2.1.26 hold. Then $f\left(g_{0}^{t}\right)=\ldots=f\left(g_{h-1}^{t}\right)$.

Proof. Clearly, it suffices to show $f\left(g_{h-2}^{t}\right)=f\left(g_{h-1}^{t}\right)$. Apply Lemma 2.1.26 to see that for all $a_{1}, \ldots, a_{h-2} \in A^{\prime},\left(a_{1}, \ldots, a_{h-2}, f\left(g_{h-2}^{t}\right), f\left(g_{h-1}^{t}\right)\right) \in D$ by choosing $b_{1}, \ldots, b_{h-2}$ from $\left\{\pi_{1}^{q}, \ldots, \pi_{q}^{q}\right\}$. Then we must have $f\left(g_{h-2}^{t}\right)=f\left(g_{h-1}^{t}\right)$ because if $\left(a_{1}, \ldots, a_{h-2}, a, b\right) \in D$ for all $a_{1}, \ldots, a_{h-2} \in A^{\prime}$, then $(a, b) \in E$; thus, $a$ and $b$ represent the same equivalence class of $E$ so that they must be equal.
2.1.28 Lemma. Let $g: h^{\kappa} \rightarrow h, t \in h^{\kappa}, b_{p} \in B_{p-1}^{t}$ for $1 \leq p \leq h$ and $\left(f\left(g_{0}^{t}\right), \ldots, f\left(g_{h-1}^{t}\right)\right) \notin D$. Consider $\left(f\left(b_{1}\right), \ldots, f\left(b_{h}\right)\right)$; then for all $1 \leq r<s \leq h$ the following holds: If we replace the $r$-th component, $f\left(b_{r}\right)$, by $f\left(g_{r-1}^{t}\right)$ and the s-th component, $f\left(b_{s}\right)$, by $f\left(g_{s-1}^{t}\right)$, then the resulting tuple is not an element of $D$.

Proof. It suffices to prove the assertion for $r=h-1$ and $s=h$. We will apply a similar induction as in the proof of Lemma 2.1.26: Consider $\left(b_{1}, \ldots, b_{n}, g_{n}^{t}, \ldots, g_{h-1}^{t}\right)$ where $n<h-2$. Then replacing any component other than $g_{n}^{t}$ by $b_{n+1}$ yields an element of $\omega_{h^{\kappa}}$; for $b_{n+1}(t)=g_{n}^{t}(t)=n$ and for $t^{\prime} \neq t, g_{n}^{t}\left(t^{\prime}\right)=g_{h-2}^{t}\left(t^{\prime}\right)=g_{h-1}^{t}\left(t^{\prime}\right)=g\left(t^{\prime}\right)$. Therefore, application of $f$ to such a tuple gives us a member of $D$ so that if our induction assumption is $\left(f\left(b_{1}\right), \ldots, f\left(b_{n+1}\right), f\left(g_{n+1}^{t}\right), \ldots, f\left(g_{h-1}^{t}\right)\right) \in D$, then the homogeneity of $D$ implies
$\left(f\left(b_{1}\right), \ldots, f\left(b_{n}\right), f\left(g_{n}^{t}\right), \ldots, f\left(g_{h-1}^{t}\right)\right) \in D$. But as we assume that $\left(f\left(g_{0}^{t}\right), \ldots, f\left(g_{h-1}^{t}\right)\right) \notin D$, by our induction we have $\left(f\left(b_{1}\right), \ldots, f\left(b_{h-2}\right), f\left(g_{h-2}^{t}\right), f\left(g_{h-1}^{t}\right)\right) \notin D$.
2.1.29 Lemma. Let $t \in h^{\kappa}$. Then either for all $g: h^{\kappa} \rightarrow h$ we have $f\left(g_{0}^{t}\right)=\ldots=f\left(g_{h-1}^{t}\right)$ or for all $g: h^{\kappa} \rightarrow h$ we have $\left(f\left(g_{0}^{t}\right), \ldots, f\left(g_{h-1}^{t}\right) \notin D\right.$.

Proof. Suppose that for some $g: h^{\kappa} \rightarrow h,\left(f\left(g_{0}^{t}\right), \ldots, f\left(g_{h-1}^{t}\right)\right) \in D$. By Lemma 2.1.27, $f\left(g_{0}^{t}\right)=\ldots=f\left(g_{h-1}^{t}\right)$. Let $\tilde{g}: h^{\kappa} \rightarrow h$. It is easy to verify $\left(\tilde{g}_{0}^{t}, \ldots, \tilde{g}_{h-2}^{t}, g_{0}^{t}\right) \in \omega_{h^{\kappa}}$, and therefore, $\left(f\left(\tilde{g}_{0}^{t}\right), \ldots, f\left(\tilde{g}_{h-2}^{t}\right), f\left(g_{0}^{t}\right)\right) \in D$. As $f\left(g_{0}^{t}\right)=f\left(g_{h-1}^{t}\right)$, this tuple can be written as $\left(f\left(\tilde{g}_{0}^{t}\right), \ldots, f\left(\tilde{g}_{h-2}^{t}\right), f\left(g_{h-1}^{t}\right)\right) \in D$. But then assuming $\left(f\left(\tilde{g}_{0}^{t}\right), \ldots, f\left(\tilde{g}_{h-1}^{t}\right)\right) \notin D$ and application of Lemma 2.1.28 by replacing the first two components of $\left(f\left(\tilde{g}_{0}^{t}\right), \ldots, f\left(\tilde{g}_{h-2}^{t}\right), f\left(g_{h-1}^{t}\right)\right)$ leads to a contradiction: The vector stays the same but is supposed to result in a vector not in $D$. Therefore, $\left(f\left(\tilde{g}_{0}^{t}\right), \ldots, f\left(\tilde{g}_{h-1}^{t}\right)\right) \in D$ so that $f\left(\tilde{g}_{0}^{t}\right)=\ldots=f\left(\tilde{g}_{h-1}^{t}\right)$.

Let $T=\left\{t_{1}, \ldots, t_{\lambda}\right\}$ be the subset of $h^{\kappa}$ containing all $t$ such that for some (or all) $g: h^{\kappa} \rightarrow h$, $\left(f\left(g_{0}^{t}\right), \ldots, f\left(g_{h-1}^{t}\right)\right) \notin D$. Denote by $S$ the complement of $T$ in $h^{\kappa}$, that is, $t \in S$ iff for some (or all) $g: h^{\kappa} \rightarrow h, f\left(g_{0}^{t}\right)=\ldots=f\left(g_{h-1}^{t}\right)$. For $g: h^{\kappa} \rightarrow h$, let $\hat{g}=\left(g\left(t_{1}\right), \ldots, g\left(t_{\lambda}\right)\right) \in h^{\lambda}$.
2.1.30 Lemma. Let $g_{1}, g_{2}: h^{\kappa} \rightarrow h$ with $\hat{g}_{1}=\hat{g}_{2}$. Then $f\left(g_{1}\right)=f\left(g_{2}\right)$.

Proof. Since $\hat{g}_{1}=\hat{g}_{2}, g_{1}$ and $g_{2}$ differ only on $S$. But if $s \in S$, we have by the previous lemma that for any $g: h^{\kappa} \rightarrow h$ and all $1 \leq i \leq h-1, f\left(g_{i}^{s}\right)=f(g)$. Thus, we may alter the values of $g_{1}$ on $S$ to those of $g_{2}$ without changing its image under $f$ and the assertion follows.
2.1.31 Lemma. Let $g_{1}, g_{2}: h^{\kappa} \rightarrow h$ with $\hat{g}_{1} \neq \hat{g}_{2}$. Then $f\left(g_{1}\right) \neq f\left(g_{2}\right)$.

Proof. Let $t \in T$ with $g_{1}(t) \neq g_{2}(t)$. Assume without loss of generality that $g_{1}(t)=0$ and $g_{2}(t)=h-1$. Then clearly, $\left(g_{1}\right)_{0}^{t}=g_{1}$ and $\left(g_{2}\right)_{h-1}^{t}=g_{2}$. By definition of $T, \quad\left(f\left(\left(g_{1}\right)_{0}^{t}\right), \ldots, f\left(\left(g_{1}\right)_{h-1}^{t}\right)\right) \notin D$ so that application of Lemma 2.1.28 by replacement of the first two components of $\left(f\left(\left(g_{2}\right)_{0}^{t}\right), \ldots, f\left(\left(g_{2}\right)_{h-1}^{t}\right)\right)$ yields $\left(f\left(g_{1}\right), f\left(\left(g_{1}\right)_{1}^{t}\right), f\left(\left(g_{2}\right)_{2}^{t}\right), \ldots, f\left(\left(g_{2}\right)_{h-2}^{t}\right), f\left(g_{2}\right)\right) \notin D$. Therefore, since $D$ is totally reflexive, $f\left(g_{1}\right) \neq f\left(g_{2}\right)$.
2.1.32 Lemma. $\left|A^{\prime}\right|=h^{\lambda}$.

Proof. Clearly, $f: h^{h^{\kappa}} \rightarrow A^{\prime}=\left\{\alpha_{1}, \ldots, \alpha_{q}\right\}$ is onto since $f\left(\pi_{j}^{\kappa}\right)=\alpha_{j}$ for $1 \leq j \leq q$. Thus, by Lemmas 2.1.30 and 2.1.31, $\left|A^{\prime}\right|=\left|\left\{f(g) \mid g \in h^{h^{\kappa}}\right\}\right|=\left|\left\{\hat{g} \mid g \in h^{h^{\kappa}}\right\}\right|=\left|h^{\lambda}\right|$, the latter equality holding because for every tuple $\left(i_{1}, \ldots, i_{\lambda}\right) \in h^{\lambda}$ there is a $g \in h^{h^{\kappa}}$ such that $g\left(t_{j}\right)=i_{j}$, $j=1, \ldots, \lambda$.

We define $\varphi^{\prime}: A^{\prime} \rightarrow h^{\lambda}$ as follows: For $a \in A^{\prime}, \varphi^{\prime}(a)=\hat{g}$, where $g$ is an arbitrary element of $h^{h^{\kappa}}$ satisfying $f(g)=a$. By Lemma 2.1.31, $\varphi^{\prime}$ is well-defined, by Lemma 2.1.30 it is one-one and so together with Lemma 2.1.32 we have that $\varphi^{\prime}$ is a bijection.
2.1.33 Lemma. Let $g_{1}, \ldots, g_{h}: h^{\kappa} \rightarrow h$ such that there is $t \in T$ with $\left(g_{1}(t), \ldots, g_{h}(t)\right) \notin \iota_{h}^{h}$. Then $\left(f\left(g_{1}\right), \ldots, f\left(g_{h}\right)\right) \notin D$.

Proof. Assume without loss of generality that $g_{i}(t)=i-1,1 \leq i \leq h$. Let $\tilde{0}: h^{\kappa} \rightarrow$ $h$ be the zero function, that is, $\tilde{0}(s)=0$ for all $s \in h^{\kappa}$. Then for all $0 \leq i, j \leq h-1$ with $i \neq j,\left(\tilde{0}_{0}^{t}, \ldots, \tilde{0}_{i-1}^{t}, g_{j+1}, \tilde{0}_{i+1}^{t}, \ldots, \tilde{0}_{h-1}^{t}\right) \in \omega_{h^{k}}$ : For $s \neq t$, evaluating the tuple at $s$ yields $\left(0, \ldots, 0, g_{j+1}(s), 0, \ldots, 0\right) \in \iota_{h}^{h}$ and evaluating the tuple at $t$ results in $(0, \ldots, i-1, j, i+$ $1, \ldots, h-1) \in \iota_{h}^{h}$ as $i \neq j$. Thus, $\left(f\left(\tilde{0}_{0}^{t}\right), \ldots, f\left(\tilde{0}_{i-1}^{t}\right), f\left(g_{j+1}\right), f\left(\tilde{0}_{i+1}^{t}\right), \ldots, f\left(\tilde{0}_{h-1}^{t}\right)\right) \in D$. Now suppose $\left(f\left(g_{1}\right), \ldots, f\left(g_{h}\right)\right) \in D$; then if we make use of the strong homogeneity of $D$ by taking $\left(f\left(g_{1}\right), \ldots, f\left(g_{h}\right)\right)$ as a replacement vector we get that $\left(f\left(\tilde{0}_{0}^{t}\right), \ldots, f\left(\tilde{0}_{h-1}^{t}\right)\right) \in D$. Hence by Lemma 2.1.27, $f\left(\tilde{0}_{0}^{t}\right)=\ldots=f\left(\tilde{0}_{h-1}^{t}\right)$. But $t \in T$ and obviously $\tilde{0}_{0}^{t}(t)=0 \neq 1=\tilde{0}_{1}^{t}(t)$ so that by Lemma 2.1.31, $f\left(\tilde{0}_{0}^{t}\right) \neq f\left(\tilde{0}_{1}^{t}\right)$, contradiction. Therefore, we must have $\left(f\left(g_{1}\right), \ldots, f\left(g_{h}\right)\right) \notin$ D.

The last step is to show
2.1.34 Lemma. $\varphi^{\prime}(D)=\omega_{\lambda}$.

Proof. Let $\left(a_{1}, \ldots, a_{h}\right) \in D$, and choose $g_{i}: h^{\kappa} \rightarrow h$ for $1 \leq i \leq h$ such that $f\left(g_{i}\right)=a_{i}$. Then by the previous lemma, $\left(g_{1}(t), \ldots, g_{h}(t)\right) \in \iota_{h}^{h}$ for all $t \in T$ so that $\left(\hat{g}_{1}, \ldots, \hat{g}_{h}\right)=\left(\varphi^{\prime}\left(a_{1}\right), \ldots, \varphi^{\prime}\left(a_{h}\right)\right) \in \omega_{\lambda}$. Conversely every element of $\omega_{\lambda}$ can clearly be written as $\left(\hat{g}_{1}, \ldots, \hat{g}_{h}\right)=\left(\varphi^{\prime}\left(a_{1}\right), \ldots, \varphi^{\prime}\left(a_{h}\right)\right)$ for some $g_{i}: h^{\kappa} \rightarrow h$ and some $a_{i} \in A^{\prime}$, $1 \leq i \leq h$. We may assume that $g_{i}(s)=0$ for all $s \in S$ and all $1 \leq i \leq h$. But then for arbitrary $u \in h^{\kappa},\left(g_{1}(u), \ldots, g_{h}(u)\right) \in \iota_{h}^{h}$ so that $\left(g_{1}, \ldots, g_{h}\right) \in \omega_{h^{\kappa}}$. Hence, $\left(f\left(g_{1}\right), \ldots, f\left(g_{h}\right)\right)=\left(a_{1}, \ldots, a_{h}\right) \in D$.
2.1.35 Theorem. $C$ is h-regularly generated.

Proof. By Lemma 2.1.21, $\gamma^{-1}(D)=C^{*}=C$ so that by the previous lemma, $C=\left(\varphi^{\prime} \circ \gamma\right)^{-1}\left(\omega_{\lambda}\right)$. Since Lemma 2.1.32 implies that $\varphi^{\prime} \circ \gamma: C \rightarrow \omega_{\lambda}$ is onto, the theorem follows.

The case $m=2$
Recall that in the beginning we established that for $m \geq 4, \mathfrak{B}$ must be totally reflexive and totally symmetric. We showed that this case is impossible; now we will consider the case $m=2$.

First note that for $0 \in \operatorname{Con}(\mathfrak{A})$ we have that either $0 \cap B=0$ (that is, $0 \subseteq B$ ) or $0 \cap B=\emptyset$; for otherwise, the projection of that intersection on one coordinate would be a proper subalgebra of $\mathfrak{A}$. In the first case $\mathfrak{B}$ is reflexive, in the second case we call $\mathfrak{B}$ areflexive. Observe that if $\mathfrak{B}$ is areflexive then, as a binary relation, it can neither have a least nor a greatest element.

For two binary relations $C_{1}, C_{2}$ we denote the relation product by $C_{1} \cdot C_{2}$; we define the inverse relation of $C_{1}$ to be $C_{1}^{-1}=\left\{(a, b) \mid(b, a) \in C_{1}\right\}$.
2.1.36 Lemma. If $B \cdot B^{-1}=A^{2}$, then $B$ has a least and a greatest element.

Proof. We will prove that $B$ has a greatest element. Since $B \cdot B^{-1}=A^{2}$ if and only if $B^{-1} \cdot B=$ $A^{2}$ this implies that $B^{-1}$ has a greatest element as well so that $B$ has a least element. Define for $2 \leq i \leq \kappa$ sets $C_{i} \subseteq A^{i}$ by

$$
C_{i}=\left\{\left(a_{1}, \ldots, a_{i}\right) \mid \exists b \in A \forall 1 \leq j \leq i\left(a_{j}, b\right) \in B\right\} .
$$

Then clearly $C_{2}=B \cdot B^{-1}=A^{2}$ and $\left(C_{i}, F\right)$ is a totally symmetric subalgebra of $\mathfrak{A}^{i}$ for all $2 \leq i \leq \kappa$. Moreover, $C_{i}=A^{i}$ obviously implies that $C_{i+1}$ is totally reflexive so that it must be equal to $A^{i+1}$. By induction we get $C_{\kappa}=A^{\kappa}$ and so $\left(\alpha_{1}, \ldots, \alpha_{\kappa}\right) \in C_{\kappa}$. But that means there exists $o \in A$ such that $(a, o) \in B$ for all $a \in A$; hence, $o$ is a greatest element of $B$.

Now if $B$ is areflexive, it has no greatest element so that by the previous lemma, $B \cdot B^{-1} \neq A^{2}$. Since $B \cdot B^{-1}$ is reflexive, $\left|B \cdot B^{-1}\right| \geq|A|$. Furthermore, $|B|>|A|$ so that there exist $a, b, c \in A$ with $a \neq b$ such that $(a, c) \in B$ and $(b, c) \in B$. Thus, $(a, b) \in B \cdot B^{-1}$ and so $\left|B \cdot B^{-1}\right|>|A|$. Hence, as $\left(B \cdot B^{-1}, F\right)$ is a subalgebra of $\mathfrak{A}$, we can consider $B \cdot B^{-1}$ instead of $B$; we will therefore assume from now on that $B$ is reflexive.

We call $B$ antisymmetric iff $B \cap B^{-1}=0$. Suppose $B$ is not antisymmetric; then since $B$ is reflexive, $B \cap B^{-1}$ properly contains 0 . In that case we may as well assume $B$ is symmetric by replacing $B$ with $B \cap B^{-1}$ which trivially is symmetric. Hence, we are back in the totally reflexive totally symmetric case which we have already shown absurd.

So we assume $B$ is antisymmetric as well. Then the following is true.
2.1.37 Lemma. B has exactly one least and exactly one greatest element.

Proof. Note that $B$ can have at most one least and one greatest element because of its antisymmetry. Since $B \cdot B^{-1}$ contains 0 properly and since it is symmetric, it must equal $A^{2}$ so that by the last lemma, $B$ has at least one least and at least one greatest element.
2.1.38 Lemma. $B \cdot B \neq A^{2}$.

Proof. Let $a \in A$ and let $o$ be the greatest element of $B$. If $(o, a) \in B \cdot B$, then there exists $b \in A$ such that $(o, b) \in B$ and $(b, a) \in B$. Because $o$ is the greatest element of $B,(b, o) \in B$ and so, since $B$ is antisymmetric, $b=o$. Hence, $(o, a) \in B$ so that $o=a$. Thus, if $a \neq o$, then $(o, a) \notin B \cdot B$.

Assume that $\mathfrak{B}$ is maximal among the antisymmetric subalgebras of $\mathfrak{A}^{2}$. We will finish the case $m=2$ and show $B$ is a partial order with least and greatest element; the only thing that is missing is the transitivity if $B$.
2.1.39 Lemma. $B$ is transitive.

Proof. Consider the subalgebra $(B \cdot B, F)$ of $\mathfrak{A}^{2}$. By the previous lemma, $B \cdot B \neq A^{2}$. Suppose $B \cdot B$ is not antisymmetric. Then $(B \cdot B) \cap(B \cdot B)^{-1}$ is strictly between 0 and $1 \in \operatorname{Con}(\mathfrak{A})$; since it is both symmetric and reflexive, this is impossible. Thus, $B \cdot B$ is antisymmetric, and as $B \cdot B \supseteq B$, the maximality of $B$ implies $B \cdot B=B$. Hence, $B$ is transitive.

We are left with the case $m=3$ as we have already eliminated the cases $m \geq 4$ and $m=2$ as possibilities.

## The case $m=3$

We will show that the case where $\mathfrak{B} \leq \mathfrak{A}^{3}$ is impossible as well to finish the proof. Denote by $\Delta_{n}(A)$ the diagonal of $A^{n}$.
2.1.40 Lemma. $B_{(1,2,3)}=\Delta_{3}(A)$, that is, $B$ contains $\Delta_{3}(A)$.

Proof. Since $B_{(1,2)}$ is essentially a subuniverse of $\mathfrak{A}^{2}$ we have that $\left|B_{(1,2)}\right|$ equals either $|A|^{2}$ or $|A|$. In the first case our assertion follows immediately. In the second case consider the subalgebra $\left(\rho_{1}\left(B_{(1,2)}\right), F\right)$ of $\mathfrak{A}^{2}$ and note that $\left|\rho_{3}(B)\right|=|A|^{2}$ implies that $\pi_{1}^{2}\left(\rho_{1}\left(B_{(1,2)}\right)\right)=A$. But since $\left|\rho_{1}\left(B_{(1,2)}\right)\right|=\left|B_{(1,2)}\right|=|A|, \rho_{1}\left(B_{(1,2)}\right)$ is the graph of an automorphism of $\mathfrak{A}$ and must therefore equal $\Delta_{2}(A)$. Hence, $B_{(1,2)}=\Delta_{3}(A)$ and so also $B_{(1,2,3)}=\Delta_{3}(A)$.

We will show now that we can assume that $B \cap \iota_{3}^{A} \supsetneqq \Delta_{3}(A)$. Define a subuniverse $B^{\prime}$ of $\mathfrak{A}^{3}$ by

$$
B^{\prime}=\left\{\left(a_{2}, a_{3}, a_{3}^{\prime}\right) \mid \exists a_{1} \in A\left(\left(a_{1}, a_{2}, a_{3}\right) \in B \wedge\left(a_{1}, a_{2}, a_{3}^{\prime}\right) \in B\right)\right\}
$$

Since $\left|\rho_{1}(B)\right|=|A|^{2}$ by Lemma 2.1.6, $\left|B_{(2,3)}^{\prime}\right|=|A|^{2}$. Therefore, $B^{\prime} \cap \iota_{3}^{A} \neq \Delta_{3}(A)$. It is possible that $B \cap \iota_{3}^{A}=\Delta_{3}(A)$. But in that case, $\left|B_{(1,2)}^{\prime}\right|=|A|$ and thus $\left|B^{\prime}\right|<|A|^{3}$; moreover, since $\left|\rho_{3}(B)\right|=|A|^{2}$ and $|B|>|A|^{2},\left|B^{\prime}\right|>|A|^{2}$. Hence, we can replace $B$ by $B^{\prime}$ in our proof and so we will assume from now on that $B \cap \iota_{3}^{A} \neq \Delta_{3}(A)$. Up to symmetry, this leaves us with three possibilities (since $B_{(i, j)}$ is essentially a subuniverse of $\mathfrak{A}^{2}$ and must therefore have cardinality a power of $|A|$ ):

### 2.1.41 Lemma. Either

1. $\left|B_{(1,2)}\right|=|A|^{2}$ and $\left|B_{(1,3)}\right|=\left|B_{(2,3)}\right|=|A|$ or
2. $\left|B_{(1,2)}\right|=\left|B_{(1,3)}\right|=|A|^{2}$ and $\left|B_{(2,3)}\right|=|A|$ or
3. $\left|B_{(1,2)}\right|=\left|B_{(1,3)}\right|=\left|B_{(2,3)}\right|=|A|^{2}$

### 2.1.42 Lemma. Possibility 3 is impossible.

Proof. If Possibility 3 was true, then $\iota_{3}^{A} \subseteq B$ and so also the subuniverse generated by $\iota_{3}^{A}$ would be a subset of $B$ and thus a proper subuniverse of $\mathfrak{A}^{3}$. But this is impossible as that subuniverse is totally reflexive and totally symmetric.
2.1.43 Lemma. Possibility 1 is impossible.

Proof. Define for $2 \leq i \leq \kappa$ sets $B_{i} \subseteq A^{i}$ by

$$
B_{i}=\left\{\left(a_{1}, \ldots, a_{i}\right) \mid \exists a \in A \exists b \in A \forall 1 \leq j \leq i \quad\left(a_{j}, a, b\right) \in B\right\}
$$

Clearly, $B_{i}$ is a totally symmetric subuniverse of $\mathfrak{A}^{i}$. Since $\Delta_{3}(A) \subseteq B, \Delta_{2}(A) \subseteq B_{2}$. Also, $|B|>|A|^{2}$ and so there are distinct $a, b, c \in A$ such that $(a, b, c) \in B$. But also $(b, b, c) \in B$ so that $(a, b) \in B_{2}$, and therefore $\left|B_{2}\right|>|A|$. Hence, as $\left|B_{2}\right|$ must be a power of $|A|$, we must have $\left|B_{2}\right|=|A|^{2}$. Now since $B_{n}=A^{n}$ implies that $B_{n+1}$ is totally reflexive, it implies further $B_{n+1}=A^{n+1}$. By induction we get $B_{\kappa}=A^{\kappa}$ and as a consequence, $\left(\alpha_{1}, \ldots, \alpha_{\kappa}\right) \in B_{\kappa}$. But that means that there exist $a, b \in A$ such that for all $x \in A$ we have $(x, a, b) \in B$. Setting $x=b$ yields $(b, a, b) \in B$, and since $\left|B_{(1,3)}\right|=|A|$ and $\Delta_{3}(A) \subseteq B$, we conclude $a=b$. But if we choose now $x \neq a$, then we get $(x, a, a) \in B$, contradicting $\left|B_{(2,3)}\right|=|A|$. Hence, Case 1 cannot occur.

We will conclude the proof of Theorem 2.1.4 by showing that Possibility 2 in 2.1.41 is impossible as well. This will require more effort than the other cases.
2.1.44 Lemma. If for $a, b, c \in A$ both $(a, b, c) \in B$ and $(b, a, c) \in B$, then $a=b$.

Proof. Define a subuniverse $B^{\prime}$ of $\mathfrak{A}^{3}$ by

$$
B^{\prime}=\{(b, a, c) \mid(a, b, c) \in B\}
$$

and consider the subuniverse $B \cap B^{\prime}$. Because $\left|B_{(2,3)}\right|=|A|,(a, b, b) \in B$ implies $a=b$ and clearly $(b, a, b) \in B^{\prime}$ implies the same; thus, $B \cap B^{\prime} \cap \iota_{3}^{A}=\{(a, a, b) \mid a, b \in A\}$. Therefore $B \cap B^{\prime}$ satisfies exactly the equalities of Case 1 , and so, if $\left|B \cap B^{\prime}\right|>|A|^{2}$, we have a contradiction. Hence, $\left|B \cap B^{\prime}\right|=|A|^{2}$, or equivalently, $B \cap B^{\prime}=B \cap B^{\prime} \cap \iota_{3}^{A}=\{(a, a, b) \mid a, b \in A\}$ which means exactly that $(a, b, c) \in B$ and $(b, a, c) \in B$ implies $a=b$.

For an equivalence relation $\sim$ on $\{1, \ldots, n\}$ we define $\Delta_{\sim} \subseteq A^{n}$ by

$$
\Delta_{\sim}=\left\{\left(a_{1}, \ldots, a_{n}\right) \mid i \sim j \rightarrow a_{i}=a_{j}\right\}
$$

In this context, we will denote an equivalence relation by writing down its equivalence classes. Let $(C, F)$ be the subalgebra of $\mathfrak{A}^{4}$ generated by

$$
\Delta_{\{1,2\}\{3,4\}} \cup \Delta_{\{1,3\}\{2,4\}} \cup \Delta_{\{1,4\}\{2,3\}} ;
$$

clearly, $C$ is totally symmetric. Also for $1 \leq i \leq 4, \rho_{i}(C) \supseteq \iota_{3}^{A}$, and so $\rho_{i}(C)=A^{3}$.
2.1.45 Lemma. For $\{i, j, r, s\}=\{1,2,3,4\}$, we have $C_{(i, j)}=C_{(i, j)(r, s)}$.

Proof. Define a subuniverse $B^{\prime}$ of $\mathfrak{A}^{4}$ by

$$
\begin{aligned}
B^{\prime}=\{ & \left(a_{1}, a_{2}, a_{3}, a_{4}\right) \mid \exists a_{5} \in A \exists a_{6} \in A \\
& \left.\left(\left(a_{1}, a_{2}, a_{5}\right) \in B \wedge\left(a_{5}, a_{3}, a_{4}\right) \in B \wedge\left(a_{2}, a_{1}, a_{6}\right) \in B \wedge\left(a_{6}, a_{3}, a_{4}\right) \in B\right)\right\}
\end{aligned}
$$

By our assumptions for Case 2 , one easily checks that $C \subseteq B^{\prime}$. Note next that if ( $a_{1}, a_{2}, a_{3}, a_{3}$ ) $\in$ $B^{\prime}$, then $a_{6}=a_{5}=a_{3}$ so that $\left(a_{1}, a_{2}, a_{3}\right) \in B$ and $\left(a_{2}, a_{1}, a_{3}\right) \in B$ which implies $a_{1}=a_{2}$. The same property holds for $C$ since $C \subseteq B^{\prime}$. Hence, $C_{(3,4)}=C_{(3,4)(1,2)}$ and the lemma follows from the total symmetry of $C$.
2.1.46 Lemma. Let $a_{1}, a_{2}, a_{3} \in A$ be given. Then there exists exactly one $a_{4} \in A$ such that $\left(a_{1}, a_{2}, a_{3}, a_{4}\right) \in C$.

Proof. The existence of such an $a_{1}$ follows immediately from $\rho_{1}(C)=A^{3}$. We will show that $a_{1}$ is unique. Define $B^{\prime} \leq A^{3}$ to be

$$
B^{\prime}=\left\{\left(a_{1}, a_{2}, a_{3}\right) \mid \exists a_{4} \in A \exists a_{5} \in A \quad\left(\left(a_{1}, a_{3}, a_{4}, a_{5}\right) \in C \wedge\left(a_{2}, a_{3}, a_{4}, a_{5}\right) \in C\right)\right\}
$$

It is obvious that $B^{\prime}$ is a subuniverse of $\mathfrak{A}^{3}$. Clearly, $\left|B_{(1,2)}^{\prime}\right|=|A|^{2}$. If $\left(a_{1}, a_{2}, a_{2}\right) \in B^{\prime}$, then by the definition of $B^{\prime},\left(a_{2}, a_{2}, a_{4}, a_{5}\right) \in C$ so that by the previous lemma, $a_{1}=a_{2}$. Thus, $\left|B_{(2,3)}^{\prime}\right|=|A|$, and by the same argument, $\left|B_{(1,3)}^{\prime}\right|=|A|$. But as we have already proven Case 1 impossible, we must have $\left|B^{\prime}\right|=|A|^{2}$ which implies $B^{\prime}=B_{(1,2)}^{\prime}$. So if we have $\left(a_{1}, a_{2}, a_{3}, a_{4}\right) \in C$ and $\left(a_{1}^{\prime}, a_{2}, a_{3}, a_{4}\right) \in C$, then $a_{1}=a_{1}^{\prime}$.

A consequence of the previous lemma is that we can define a function $f: A^{3} \rightarrow A$ assigning to $\left(a_{2}, a_{3}, a_{4}\right) \in A^{3}$ the unique $a_{1} \in A$ such that $\left(a_{1}, a_{2}, a_{3}, a_{4}\right) \in C$.
2.1.47 Lemma. For all permutations $\pi$ of $\{a, b, c\}$ we have that $f(a, b, c)=f(\pi(a), \pi(b), \pi(c))$. Moreover, $f(a, a, b)=b$, and $f(a, b, f(a, b, c))=c$.

Proof. The first assertion holds because $C$ is totally symmetric. The second one is a direct consequence of Lemma 2.1.45. For the last one, note that by the definition of $f,(a, b, c, f(a, b, c)) \in C$. Thus, again by the definition of $f$ and the total symmetry of $C, f(a, b, f(a, b, c))=c$.
2.1.48 Lemma. $f$ satisfies the equation $f(a, b, c)=f(f(a, d, c), d, b)$.

Proof. Define $C^{\prime}$ to be the subset of $A^{4}$ containing exactly the tuples $(a, b, c, d)$ for which $f$ satisfies the equation of the lemma; $C^{\prime}$ is a subuniverse of $\mathfrak{A}^{4}$. It is easy to check with the properties of $f$ stated in Lemma 2.1.47 that for all $1 \leq i<j \leq 4,\left|C_{(i, j)}^{\prime}\right|=|A|^{3}$. Thus, $\left|C^{\prime}\right|>|A|^{3}$ and consequently $C^{\prime}=A^{4}$. This proves the lemma.

Now choose $0 \in A$ arbitrarily and define a binary operation + on $A$ by

$$
a+b=f(a, b, 0)
$$

This will give us a prime affine relation and lead the last possibility ad absurdum.
2.1.49 Lemma. $(A,+)$ is an abelian 2-group.

Proof. Checking the associative law, we use Lemmas 2.1.47 and 2.1.48 to calculate

$$
\begin{aligned}
a+(b+c) & =f(a, b+c, 0) \\
& =f(a, f(b, c, 0), 0) \\
& =f(f(b, 0, c), 0, a) \\
& =f(b, a, c) \\
& =f(a, b, c)
\end{aligned}
$$

A similar computation yields $(a+b)+c=f(a, b, c)$ so that $a+(b+c)=(a+b)+c .0 \in A$ is the neutral element since for all $a \in A, a+0=0+a=f(a, 0,0)=a$ by Lemma 2.1.47. As by Lemma 2.1.47 we have $a+a=f(a, a, 0)=0$, each element is its own inverse. Observe that this also means the group $(A,+)$ is a 2-group. Finally, the group is abelian since $a+b=f(a, b, 0)=$ $f(b, a, 0)=b+a$ for all $a, b \in A$.

Now all tuples in $C$ are of the form $(a, b, c, f(a, b, c)$ ), which we know can be written as $(a, b, c, a+b+c)$. If we set $c=0$ then by the fact that $C$ is a subuniverse of $\mathfrak{A}^{4}$ we get that for any $n$-ary operation $g$ of $\mathfrak{A}, g(a)+g(b)+g(0)=g(a+b)$, where $a, b \in A^{n}$ are arbitrary. Define $\rho \subseteq A^{4}$ by

$$
(a, b, c, d) \in \rho \leftrightarrow a+b=c+d
$$

Then $\rho$ is a subuniverse of $\mathfrak{A}^{4}$ since for $n$-ary $g$ we have that if $a+b=c+d$, where $a, b, c, d \in A^{n}$, then $g(a+b)=g(c+d)$ and so by the preceding discussion $g(a)+g(b)=g(c)+g(d)$. Hence, $\rho$ is a prime affine relation with respect to the abelian 2 -group $(A,+)$ and so forbidden by Rosenberg's list. We have therefore eliminated Possibility 2 as a possibility in 2.1.41 and so finally finished the case $m=3$. Theorem 2.1.4 has been proven.

## $2.2 \mathcal{V}(\mathfrak{A})$ is congruence permutable

We will use the result of the last section, namely that $\mathfrak{A}$ has almost minimal spectrum, to show that the variety generated by $\mathfrak{A}$ is congruence permutable. In the beginning of our proof, we will follow another result by R. Quackenbush in [9]; for the second part we will go another way than the one shown there.
2.2.1 Definition. An algebra $\mathfrak{A}$ is called congruence permutable iff for all congruences $\psi, \theta$ on $\mathfrak{A}, \psi \cdot \theta=\theta \cdot \psi$. We say a variety is congruence permutable iff every algebra in the variety is.
2.2.2 Theorem (R. Quackenbush [9]). Let $\mathfrak{A}$ be a finite non-trivial algebra. If $\mathfrak{A}$ has minimal spectrum, then the variety generated by $\mathfrak{A}$ is congruence permutable.
2.2.3 Remark. The converse holds under the assumption that $\mathfrak{A}$ is simple and has no proper subalgebras. For a proof of this consult [9]. Note that the assumption of the theorem is that $\mathfrak{A}$ has minimal spectrum, whereas we only know until now that $\mathfrak{A}$ has almost minimal spectrum. For the proof of congruence permutability, this is still sufficient.
2.2.4 Theorem. Let $\mathfrak{A}$ be a finite non-trivial algebra. If $\mathfrak{A}$ has almost minimal spectrum, then the variety generated by $\mathfrak{A}$ is congruence permutable.

For a subdirect product $\mathfrak{B}$ of algebras $\mathfrak{A}_{1}, \ldots, \mathfrak{A}_{n}$, the meet of the kernels of all projections onto the $\mathfrak{A}_{i}$ is clearly the trivial congruence 0 . Recall that $\mathfrak{B}$ is called an irreducible subdirect product iff for every proper subset of the projections the meet of the kernels of the projections of that subset is strictly greater than 0 .
2.2.5 Definition. A set $\mathcal{A}$ of finite algebras is a direct factor set iff whenever $\mathfrak{B}$ is an irreducible subdirect product of algebras $\mathfrak{A}_{1}, \ldots, \mathfrak{A}_{n}$ of $\mathcal{A}$, then $\mathfrak{B}=\mathfrak{A}_{1} \times \ldots \times \mathfrak{A}_{n}$.

If we look at a congruence relation $\theta$ on an algebra $\mathfrak{A}$ as a subalgebra of $\mathfrak{A}^{2}$, we will denote this subalgebra by $\mathfrak{A}_{\theta}$. Note that $\mathfrak{A}_{\theta}$ is a subdirect power of $\mathfrak{A}$.
2.2.6 Lemma. All algebras in a direct factor set are simple.

Proof. Let $\theta$ be a congruence relation on a member $\mathfrak{A}$ of a direct factor set. If $\mathfrak{A}_{\theta}$ is reducible, then its projection onto one coordinate is one-one and hence $\theta=0$. On the other hand, if it is irreducible then by the definition of a direct factor set $\theta=1$.
2.2.7 Theorem. If an algebra $\mathfrak{A}$ has almost minimal spectrum then $\{\mathfrak{A}\}$ is a direct factor set.

Proof. We must prove that an irreducible subdirect product $\mathfrak{B}$ of $n$ copies of $\mathfrak{A}$ is equal to $\mathfrak{A}^{n}$. The proof will be by induction. For $n=1$ the assertion is trivial. Assume it is true for $n-1$ and let $\mathfrak{B}$ be an irreducible subdirect product of $n$ copies of $\mathfrak{A}$. Set $\mathfrak{B}^{\prime}=\mathfrak{B} / \operatorname{ker}\left(\pi_{n}^{n}\right)$, where $\pi_{n}^{n}$ denotes the projection onto the $n$-th coordinate. As one can easily see, $\mathfrak{B}^{\prime}$ is essentially an irreducible subdirect product of $n-1$ copies of $\mathfrak{A}$ and hence by the induction assumption, $\mathfrak{B}^{\prime} \cong \mathfrak{A}^{n-1}$. As $\mathfrak{B}$ is also irreducible, $|\mathfrak{B}|>\left|\mathfrak{B}^{\prime}\right|=\left|\mathfrak{A}^{n-1}\right|$. But since $\mathfrak{B} \subseteq \mathfrak{A}^{n}$, we have that $|\mathfrak{B}| \leq\left|\mathfrak{A}^{n}\right|$ so that by the fact that $\mathfrak{A}$ has almost minimal spectrum we have $|\mathfrak{B}|=\left|\mathfrak{A}^{n}\right|$ and therefore $\mathfrak{B}=\mathfrak{A}^{n}$.
2.2.8 Definition. A congruence is uniform iff all its equivalence classes are of the same cardinality. An algebra is said to have uniform congruences iff all its congruences are uniform and a variety has uniform congruences iff all its algebras do.

Our next goal is to prove that our algebra has uniform congruences. We will need the following theorem.
2.2.9 Theorem. Let $\mathcal{A}$ be a direct factor set, $I=\{1, \ldots, n\}$ be a finite index set and $\left(\mathfrak{A}_{i}\right)_{i \in I}$ be algebras in $\mathcal{A}$. Let $\mathfrak{A}=\prod_{i \in I} \mathfrak{A}_{i}$ and let $\theta$ be a non-trivial congruence of $\mathfrak{A}$. Then there exists a proper subset $J$ of $I$ such that $\mathfrak{A}_{\theta} \cong \mathfrak{A} \times \prod_{j \in J} \mathfrak{A}_{j}$. In this case $\theta$ has $\left|\prod_{i \in I \backslash J} \mathfrak{A}_{i}\right|$ equivalence classes each of which has $\left|\prod_{j \in J} \mathfrak{A}_{j}\right|$ elements.

Proof. Since $\mathfrak{A}_{\theta} \leq \mathfrak{A}^{2}=\prod_{i \in I} \mathfrak{A}_{i} \times \prod_{i \in I} \mathfrak{A}_{i}$, we know that $\mathfrak{A}_{\theta}$ is an irreducible subdirect product and hence, since $\mathcal{A}$ is a direct factor set, a direct product of some of the factors of $\mathfrak{A}^{2}$; that is, there are subsets $K$ and $K^{\prime}$ of $I$ such that $\mathfrak{A}_{\theta} \cong \prod_{i \in K} \mathfrak{A}_{i} \times \prod_{i \in K^{\prime}} \mathfrak{A}_{i}$. We claim that each $\mathfrak{A}_{i}$ occurs at least once in this direct product representation so that $K \cup K^{\prime}=I$. Given $a_{i} \in \mathfrak{A}_{i}$, $2 \leq i \leq n$, we have that $\left(a, a_{2}, \ldots, a_{n}, a, a_{2}, \ldots, a_{n}\right) \in \mathfrak{A}_{\theta}$ for every $a \in A$ because $\theta$ is reflexive. On the other hand the components of an element in $\mathfrak{A}_{\theta}$ corresponding to indices in $K$ and $K^{\prime}$ uniquely determine the other components. Therefore, $i=1$ must be in $K$ or $K^{\prime}$ and clearly the same holds for any $i \in I$. Hence, by reordering the factors of the direct product representation of $\mathfrak{A}_{\theta}$ and setting $J=K \cap K^{\prime}$, the first assertion of the theorem follows.
Now denote the equivalence classes of $\theta$ by $C_{1}, \ldots, C_{l}$ and their cardinalities by $c_{1}, \ldots, c_{l}$. First we will show that $c_{k} \leq\left|\prod_{j \in J} \mathfrak{A}_{j}\right|, 1 \leq k \leq l$. Consider the projections

$$
\begin{array}{cccc}
\varsigma_{k}: & C_{k} & \rightarrow & \prod_{j \in J} A_{j} \\
\left(a_{1}, \ldots, a_{n}\right) & \mapsto & \left(a_{j}\right)_{j \in J}
\end{array}
$$

We claim that $\varsigma_{k}$ is one-one. Let $a, b \in C_{k}$ with $\varsigma_{k}(a)=\varsigma_{k}(b)$, that is, $\left(a_{j}\right)_{j \in J}=\left(b_{j}\right)_{j \in J}$; if we prove them equal our assertion follows. Define a vector $d$ by

$$
d_{j}= \begin{cases}b_{j} & , j \in K^{\prime} \\ a_{j} & , j \in I \backslash K^{\prime}\end{cases}
$$

As the $b_{j}$ and $a_{j}$ agree on $J$, our projection

$$
\begin{array}{cccc}
\pi: & A \times A & \rightarrow & \prod_{j \in K} A_{j} \times \prod_{j \in K^{\prime}} A_{j} \\
\left(a_{1}, \ldots, a_{n}, a_{1}^{\prime}, \ldots, a_{n}^{\prime}\right) & \mapsto & \left(\left(a_{j}\right)_{j \in K},\left(a_{j}^{\prime}\right)_{j \in K^{\prime}}\right)
\end{array}
$$

maps $(d, d) \in \theta$ to $\left(\left(a_{j}\right)_{j \in K},\left(b_{j}\right)_{j \in K^{\prime}}\right)$. But $(a, b) \in \theta$ is mapped to exactly the same vector and so, as the coordinates in $K$ and $K^{\prime}$ uniquely determine all the others, we have that $a_{j}=b_{j}$ for $j \in I \backslash K^{\prime}$. By symmetry we conclude that $a_{j}=b_{j}$ for $j \in I \backslash K$ and hence, $a=b$ follows. Finally, our inequality together with the obvious equalities $\sum_{k=1}^{l} c_{k}=|\mathfrak{A}|$ and $\sum_{k=1}^{l} c_{k}^{2}=$ $|\mathfrak{A}|\left|\prod_{j \in J} \mathfrak{A}_{j}\right|$ implies that for all $1 \leq k \leq l$ we must have $c_{k}=\left|\prod_{j \in J} \mathfrak{A}_{j}\right|$.

Now we can establish that the finite algebras in $\mathcal{V}(\mathfrak{A})$ have uniform congruences. Recall that a variety is called locally finite iff every finitely generated algebra in it is finite. For a set $\mathcal{A}$ of algebras of the same type, denote by $P(\mathcal{A})$ all products, by $S(\mathcal{A})$ all subalgebras, and by $H(\mathcal{A})$ all homomorphic images of algebras of $\mathcal{A}$. Then it is well-known that $\mathcal{V}(\mathcal{A})=H S P(\mathcal{A})$.
2.2.10 Theorem. Let $\mathcal{A}$ be a finite direct factor set with the property that $S(\mathcal{A}) \subseteq P(\mathcal{A})$. Then the finite algebras in $\mathcal{V}(\mathcal{A})$ have uniform congruences.

Proof. First let $\mathfrak{B} \in S P(\mathcal{A})$ be finite. If $\mathfrak{B}$ is a subdirect product of algebras in $\mathcal{A}$ then $\mathfrak{B} \in P(\mathcal{A})$ because $\mathcal{A}$ is a direct factor set. If on the other hand the projection $\pi_{i}$ of $B$ onto some coordinate $i$ is not onto, then by the assumption $S(\mathcal{A}) \subseteq P(\mathcal{A})$ we can replace that coordinate with a product of algebras in $\mathcal{A}$ equal to $\pi_{i}(B)$ and we have the first case again. Hence, $S P(\mathcal{A})=P(\mathcal{A})$ and thus by the last theorem, all finite algebras in $S P(\mathcal{A})$ have uniform congruences. Now let $\mathfrak{C} \in \mathcal{V}(\mathcal{A})=\operatorname{HSP}(\mathcal{A})$ be given and assume it is finite; let $\psi$ be a congruence on $\mathfrak{C}$. We want to show $\psi$ is uniform. Clearly, $\mathfrak{C} \cong \mathfrak{B} / \theta$ for some $\mathfrak{B} \in S P(\mathcal{A})$ and some (uniform) $\theta \in \operatorname{Con}(\mathfrak{B})$. We can assume $\mathfrak{B}$ is finite: Observe first that $|\mathcal{A}|<\aleph_{0}$ implies that $\mathcal{V}(\mathcal{A})$ is locally finite. Now if $\mathfrak{B}$ is infinite, replace it by the subalgebra $\mathfrak{D}$ generated by any finite subset of $B$ containing at least one representative from each $\theta$-class; then obviously $\mathfrak{C} \cong \mathfrak{D} / \tilde{\theta}$ if we set $\tilde{\theta}=\theta \cap D^{2}$. Now $\psi$ induces a congruence relation $\zeta$ on $\mathfrak{B}$, defined by $a \zeta b \leftrightarrow[a]_{\theta} \psi[b]_{\theta}$, and every congruence class of $\zeta$ corresponds to exactly one congruence class of $\psi$. Since $\zeta$ and $\theta$ are uniform, $\psi$ must be uniform as well: If the size of all $\zeta$-classes is $n$ and the size of all $\theta$-classes is $j$, then the size of all $\psi$-classes must be $\frac{n}{j}$. Hence $\psi$ is uniform.

Finally we have also established what we wanted earlier: $\mathfrak{A}$ has really minimal spectrum. Still, it is worth mentioning.
2.2.11 Theorem. Let $\mathfrak{A}$ satisfy the hypotheses of Theorem 2.0.7. Then $\mathfrak{A}$ has minimal spectrum.

Proof. This follows from the proof of the previous lemma and Theorem 2.2.9.
Our next goal is to show that $\mathcal{V}(\mathfrak{A})$ has coherent congruences.
2.2.12 Definition. An algebra $\mathfrak{A}$ is congruence coherent iff for every subalgebra $\mathfrak{B}$ of $\mathfrak{A}$ and every congruence $\theta$ on $\mathfrak{A}$ it is true that if $\mathfrak{B}$ contains a congruence class of $\theta$, then $\mathfrak{B}$ is the union of congruence classes of $\theta$; in other words, if $[a]_{\theta} \subseteq B$ for some $a \in B$ implies $[b]_{\theta} \subseteq B$ for every $b \in B$. A variety is called congruence coherent or simply coherent iff all of its members are.

The following two lemmas are due to M. Clark and P. Krauss [5].
2.2.13 Lemma. If the finite algebras in a variety $\mathcal{V}$ are congruence uniform then the finite algebras in $\mathcal{V}$ are congruence coherent.

Proof. Let $\mathfrak{A} \in \mathcal{V}$ be finite, and let $\mathfrak{B}$ be a subalgebra and $\theta$ be a congruence of $\mathfrak{A}$. If $X \subseteq B$ is an equivalence class of $\theta$, then it is also an equivalence class of $\theta \cap B^{2}$. Now let $Y$ be another congruence class of $\theta \cap B^{2}$. Then there exists a congruence class $Z$ of $\theta$ such that $Y=Z \cap B$. But by our hypothesis, $|Y|=|X|=|Z|$ and hence, since $Z$ is finite, $Y=Z$.
2.2.14 Lemma. If $\mathcal{V}$ is a locally finite variety and the finite algebras in $\mathcal{V}$ are congruence coherent then $\mathcal{V}$ is congruence coherent.

Proof. Let $\mathfrak{A} \in \mathcal{V}$ and let $\mathfrak{B}$ be a subalgebra and $\theta$ be a congruence of $\mathfrak{A}$. Consider an equivalence class $X$ of $\theta, X \subseteq B$, and consider $[b]_{\theta}$ for some $b \in B$. Now if $a \theta b$ is given, choose $x \in X$ and consider the restrictions to the subuniverse $[a, b, x]$ of $\mathfrak{A}$ generated by $\{a, b, x\}: X \cap[a, b, x]$ is a congruence class of $\theta \cap[a, b, x]$ on $\mathfrak{A} \cap[a, b, x]$ and $X \cap[a, b, x]=X \cap[a, b, x] \cap B$. Clearly, $a \equiv b(\theta \cap[a, b, x])$ and so by the hypothesis $a \in B \cap[a, b, x]$. Thus we have that $[b]_{\theta} \subseteq B$ and the lemma follows.
2.2.15 Corollary. If $\mathcal{V}$ is a locally finite variety and the finite algebras in $\mathcal{V}$ are congruence uniform then $\mathcal{V}$ is congruence coherent.

Proof. This is an immediate consequence of Lemmas 2.2.13 and 2.2.14.
We will use a version of a result on $g$-coherence from [4] to finish our proof.
2.2.16 Lemma. If a variety $\mathcal{V}$ is coherent, then for some $n$ there exist ternary terms $t_{1}, \ldots, t_{n}$ and an $n+1$-ary term $\omega$ such that the following identities hold in $\mathcal{V}$ :

$$
\begin{aligned}
& t_{i}(x, x, z)=z, \quad i=1, \ldots, n \\
& y=\omega\left(x, t_{1}(x, y, z), \ldots, t_{n}(x, y, z)\right)
\end{aligned}
$$

Proof. Consider the free algebra with three generators determined by $\mathcal{V}, \mathfrak{F}_{3}(\mathcal{V})$, and call the generators $x, y, z$. Let $\theta=\theta(x, y)$ be the congruence on $\mathfrak{F}_{3}(\mathcal{V})$ generated by $\{(x, y)\}$ and let $\mathfrak{B}$ be the subalgebra of $\mathfrak{F}_{3}(\mathcal{V})$ generated by the set $\{x\} \cup[z]_{\theta}$. Clearly, $[z]_{\theta} \subseteq B$ and hence by the coherence of $\mathcal{V},[b]_{\theta} \subseteq B$ for all $b \in B$. Since $x \in B$ and $y \theta x$, also $y \in B$ and so there exists a term $\omega$ in the language of $\mathcal{V}$ such that $y=\omega\left(x, c_{1}, \ldots, c_{n}\right)$, where $c_{1}, \ldots, c_{n} \in[z]_{\theta}$. As elements of $\mathfrak{F}_{3}(\mathcal{V})$ the $c_{i}$ have representations as terms $t_{i}(x, y, z)$ so that $y=\omega\left(x, t_{1}(x, y, z), \ldots, t_{n}(x, y, z)\right)$. Furthermore, since $t_{i}(x, y, z) \in[z]_{\theta}$ and since $\theta$ is the congruence identifying $x$ and $y$, one can easily derive that $t_{i}(x, x, z)=z$ is an identity of $\mathcal{V}$ for $i=1, \ldots, n$.

The following theorem is a well-known criterion for congruence permutability by A. Mal'cev.
2.2.17 Theorem. A variety $\mathcal{V}$ is congruence permutable iff there exists a ternary term $p(x, y, z)$ of $\mathcal{V}$ such that the identities

$$
p(x, x, z)=p(z, x, x)=z
$$

can be derived in $\mathcal{V}$. The term $p$ is called a Mal'cev term.
Proof. First assume that $\mathcal{V}$ is congruence permutable. Consider the congruences $\theta=\theta(x, y)$ and $\psi=\psi(y, z)$ generated by $\{(x, y)\}$ and $\{(y, z)\}$, respectively, on the free algebra $\mathfrak{F}_{3}(\mathcal{V})$ with generators $x, y, z$. Clearly, $(x, z) \in \theta \cdot \psi$ and so by our assumption $(x, z) \in \psi \cdot \theta$. Hence, there is a term $p(x, y, z)$ in $\mathfrak{F}_{3}(\mathcal{V})$ such that $(x, p(x, y, z)) \in \psi$ and $(p(x, y, z), z) \in \theta$ which by the definition of those congruences yields $x=p(x, z, z)$ and $p(x, x, z)=z$.

Conversely, let $p(x, y, z)$ be a term of $\mathcal{V}$ satisfying $p(x, x, z)=p(z, x, x)=z$, let $\mathfrak{A}$ be an arbitrary algebra of $\mathcal{V}$ and let $\theta, \psi \in \operatorname{Con}(\mathfrak{A})$ be any two congruences on $\mathfrak{A}$. If $(a, b) \in \theta \cdot \psi$
then there exists $c \in A$ with $(a, c) \in \theta$ and $(c, b) \in \psi$. Since trivially $(a, a),(b, b) \in \theta$ (and in $\psi),(p(a, c, b), p(a, a, b))=(p(a, c, b), b) \in \theta$ and $(p(a, b, b), p(a, c, b))=(a, p(a, c, b)) \in \psi$ proving $(a, b) \in \psi \cdot \theta$.
2.2.18 Theorem. If a variety $\mathcal{V}$ is coherent, then its congruences permute.

Proof. Set $p(x, y, z)=\omega\left(z, t_{1}(y, x, z), \ldots, t_{n}(y, x, z)\right)$, where $\omega$ and $t_{1}, \ldots, t_{n}$ are the term operations of Lemma 2.2.16. Then we have

$$
p(x, z, z)=\omega\left(z, t_{1}(z, x, z), \ldots, t_{n}(z, x, z)\right)=x
$$

and

$$
\begin{aligned}
p(x, x, z) & =\omega\left(z, t_{1}(x, x, z), \ldots, t_{n}(x, x, z)\right) \\
& =\omega(z, z, \ldots, z) \\
& =\omega\left(z, t_{1}(z, z, z), \ldots, t_{n}(z, z, z)\right) \\
& =z
\end{aligned}
$$

Hence $p(x, y, z)$ is a Mal'cev term of $\mathcal{V}$ and so by the last theorem $\mathcal{V}$ is congruence permutable.

Now we can prove this section's main theorem.
Proof of Theorem 2.2.4. Let $\mathfrak{A}$ have almost minimal spectrum. Then Lemma 2.2.7 says that $\{\mathfrak{A}\}$ is a direct factor set. Theorem 2.2.10 implies that the finite algebras in $\mathcal{V}(\mathfrak{A})$ have uniform congruences and so by Corollary 2.2.15, $\mathcal{V}(\mathfrak{A})$ is coherent. Finally reference to Theorem 2.2.18 concludes the proof.

## $2.3 \mathcal{V}(\mathfrak{A})$ is congruence distributive

In this section we will show that if $\mathfrak{A} \times \mathfrak{A}$ has a skew congruence, then $\mathfrak{A}$ is prime affine, and if not, then $\mathfrak{A}$ generates a congruence distributive equational class.
2.3.1 Definition. An algebra is called congruence distributive iff it has a distributive congruence lattice. We say a variety is congruence distributive iff all of its members are.

For algebras $\left(\mathfrak{A}_{i}\right)_{i \in I}$ of the same type there is a natural embedding of the product of their congruence lattices to the congruence lattice of their product, namely

$$
\epsilon: \begin{array}{clc}
\prod_{i \in I} \operatorname{Con}\left(\mathfrak{A}_{i}\right) & \rightarrow & \operatorname{Con}\left(\prod_{i \in I} \mathfrak{A}_{i}\right) \\
\prod_{i \in I} \theta_{i} & \mapsto & \left\{(a, b) \mid \forall i \in I\left(\left(a_{i}, b_{i}\right) \in \theta_{i}\right)\right\}
\end{array}
$$

2.3.2 Definition. A congruence on a product of algebras of the same type is called factor congruence iff it is the product of congruences on those algebras as defined before; otherwise, it is called skew.

We will now concentrate on the case where $\mathfrak{A} \times \mathfrak{A}$ has a skew congruence $\theta$; we will follow a result of $H$. P. Gumm in [8] to show that in this case $\mathfrak{A}$ is affine with respect to an abelian $p$-group for some prime $p$.
2.3.3 Definition. A lattice $\mathfrak{L}$ is called modular iff it satisfies the equation

$$
x \cap((x \cap y) \cup z)=(x \cap y) \cup(x \cap z)
$$

2.3.4 Remark. It is easy to check that a lattice $\mathfrak{L}$ is modular iff in $\mathfrak{L}, y \leq x$ implies $x \cap(y \cup z)=$ $y \cup(x \cap z)$. The 5 -element lattice $\mathcal{N}_{5}$ (over the set $\{0, a, b, c, 1\}$ it is defined by $0 \leq a \leq b \leq 1$ and $0 \leq c \leq 1$ and no other elements are comparable) is nonmodular:

$$
b \cap(a \cup c)=b \neq a=a \cup(b \cap c)
$$

Hence, every lattice containing $\mathcal{N}_{5}$ is nonmodular. Conversely, every nonmodular lattice contains a sublattice isomorphic to $\mathcal{N}_{5}$ : For if $x, y, z$ do not satisfy the modular law, then it is easily verified that the identification $(0, a, b, c, 1)=(x \cap z, y \cup(x \cap z), x \cap(y \cup z), z, y \cup z)$ is such an isomorphism. Therefore, if we have a modular lattice $\mathfrak{L}$ and two arbitrary elements $b, t \in L$ with $b \leq t$, then the length of every path from $b$ to $t$ in the Hasse diagram of $\mathfrak{L}$ is the same. In that light, the following definition makes sense.
2.3.5 Definition. For a cardinal $\alpha$, by $\mathcal{M}_{\alpha}$ we understand the modular lattice with least and greatest element and $\alpha$ atoms and no other elements.

The reason why we defined all this is the following:
2.3.6 Lemma. If $\mathfrak{A}$ is an algebra with permuting congruences, then Con $(\mathfrak{A})$ is modular.

For the proof of the lemma as well as for later proofs, we need to recall the following wellknown fact.
2.3.7 Lemma. An algebra $\mathfrak{A}$ has permuting congruences iff for all $\psi, \theta \in \operatorname{Con}(\mathfrak{A}), \psi \cup \theta=\psi \cdot \theta$.

Proof. Clearly, $\psi \cdot \theta$ is a congruence on $\mathfrak{A}$, the symmetry provided by the permutability of $\psi$ and $\theta$. If $(a, b) \in \psi$, then since trivially $(b, b) \in \theta$ we have that $(a, b) \in \psi \cdot \theta$. Hence, $\psi \cdot \theta \geq \psi$ and by the same argument $\psi \cdot \theta \geq \theta$. If $(a, b) \in \psi \cdot \theta$, then there exists $c \in A$ such that $(a, c) \in \psi$ and $(c, b) \in \theta$. Therefore any congruence $\vartheta$ with $\psi \leq \vartheta$ and $\theta \leq \vartheta$ must by its transitivity contain $(a, b)$ so that $\psi \cdot \theta \leq \vartheta$. This concludes the proof of one direction; the other one is obvious.

Proof of Lemma 2.3.6. Let $\theta_{1}, \theta_{2}, \theta_{3} \in \operatorname{Con}(\mathfrak{A})$ with $\theta_{1} \leq \theta_{2}$. It must be shown that $\theta_{2} \cap\left(\theta_{1} \cup\right.$ $\left.\theta_{3}\right)=\theta_{1} \cup\left(\theta_{2} \cap \theta_{3}\right)$, or equivalently, that $\theta_{2} \cap\left(\theta_{1} \cup \theta_{3}\right) \leq \theta_{1} \cup\left(\theta_{2} \cap \theta_{3}\right)$. Let $(a, c) \in \theta_{2} \cap\left(\theta_{1} \cup \theta_{3}\right)$; since $(a, c) \in\left(\theta_{1} \cup \theta_{3}\right)$ and since $\mathfrak{A}$ has permuting congruences, there exists $b \in A$ such that $(a, b) \in \theta_{1}$ and $(b, c) \in \theta_{3}$. Moreover, $(a, b) \in \theta_{2}$ as $\theta_{1} \leq \theta_{2}$, and since also $(a, c) \in \theta_{2}$, we have that $(b, c) \in \theta_{2} \cdot \theta_{2}=\theta_{2}$. Thus, $(b, c) \in \theta_{2} \cap \theta_{3}$; consequently, $(a, c) \in \theta_{1} \cdot\left(\theta_{2} \cap \theta_{3}\right)=\theta_{1} \cup\left(\theta_{2} \cap \theta_{3}\right)$.

Let us return to our algebra, which we know now has a modular congruence lattice. Since $\mathfrak{A}$ is simple and since $\operatorname{Con}\left(\mathfrak{A}^{2}\right) / \operatorname{ker}\left(\pi_{1}\right) \cong \operatorname{Con}\left(\mathfrak{A}^{2}\right) / \operatorname{ker}\left(\pi_{2}\right) \cong \operatorname{Con}(\mathfrak{A})$, the intervals $\left[k e r \pi_{i}, 1\right]$ are equal to $\left\{\operatorname{ker} \pi_{i}, 1\right\}$. As $\operatorname{Con}(\mathfrak{A} \times \mathfrak{A})$ is modular we conclude that the sublattice generated by $\left\{\operatorname{ker} \pi_{1}, \operatorname{ker} \pi_{2}, \theta\right\}$ is isomorphic to $\mathcal{M}_{3}$. Furthermore, the greatest (resp. least) element in $\mathcal{M}_{3}$ coincides with the greatest (resp. least) element in $\operatorname{Con}(\mathfrak{A} \times \mathfrak{A})$. We say that $\mathcal{M}_{3}$ is a 0-1-sublattice of $\operatorname{Con}(\mathfrak{A} \times \mathfrak{A})$. We summarize: $\mathfrak{A} \times \mathfrak{A}$ has three congruences $\theta_{1}, \theta_{2}, \theta_{2}$ satisfying $\theta_{i} \cdot \theta_{j}=1$ and $\theta_{i} \cap \theta_{j}=0$ for $1 \leq i, j \leq 3$. In the following, we will investigate an abstraction of this situation.

Let $S$ be a set, $|S| \geq 4$, and let $\theta_{1}, \theta_{2}, \theta_{3}$ be equivalence relations on $S$ satisfying $\theta_{i} \cdot \theta_{j}=1$ and $\theta_{i} \cap \theta_{j}=0$ for $1 \leq i, j \leq 3, i \neq j$. Then we call the quadruple $\mathcal{S}=\left(S, \theta_{1}, \theta_{1}, \theta_{1}\right)$ an $S$-3-System. A geometrical interpretation of an S-3-system, the so-called $\ddot{\text { Alquivalenzklassengeometrie, will }}$ prove useful: Call the elements of $S$ points and the equivalence classes of the relations lines. Two lines are parallel iff they are classes of the same equivalence relation. A point lies on a line iff it is an element of the line. With these definitions we have:
2.3.8 Lemma. The Äquivalenzklassengeometrie of an S-3-system has the following properties:
(S1) There are three classes of parallel lines.
(S2) Each point lies on exactly one line of each parallel-class.
(S3) Two non-parallel lines intersect in exactly one point, that is, they have exactly one point in common.

Proof. (S1) and (S2) are trivial. For (S3), let $l_{1}, l_{2}$ be two non-parallel lines, and assume without loss of generality they are equivalence classes of $\theta_{1}$ and $\theta_{2}$, respectively. Let $x \in l_{1}$ and $y \in l_{2}$ be arbitrary points. Then since $\theta_{1} \cdot \theta_{2}=1$, there is $z \in S$ such that $x \theta_{1} z$ and $z \theta_{2} y$. Hence, $z \in l_{1} \cap l_{2}$. Suppose there is another $u \in l_{1} \cap l_{2}$. Then $u \theta_{1} z$ and $u \theta_{2} z$ so that $u=z$ since $\theta_{1} \cap \theta_{2}=0$.
2.3.9 Definition. An algebra $\mathfrak{Q}=(Q, \cdot)$ with one binary operation $\cdot$ is called a quasigroup iff for all $a, c \in Q$ the equations $c \cdot x=a$ and $y \cdot c=a$ have unique solutions $x, y \in Q$.

We will show now that quasigroups give rise to S -3-systems, and conversely, from S-3-systems we can construct quasigroups. Let $\mathfrak{Q}=(Q, \cdot)$ be a quasigroup. Set $S=Q \times Q$ and define $\theta_{1}$, $\theta_{2}$ and $\theta_{2}$ on $S$ by

$$
\begin{array}{ccc}
(x, y) \theta_{1}\left(x^{\prime}, y^{\prime}\right) & \leftrightarrow & x=x^{\prime} \\
(x, y) \theta_{2}\left(x^{\prime}, y^{\prime}\right) & \leftrightarrow & y=y^{\prime} \\
(x, y) \theta_{3}\left(x^{\prime}, y^{\prime}\right) & \leftrightarrow & x \cdot y=x^{\prime} \cdot y^{\prime}
\end{array}
$$

2.3.10 Lemma. $\left(S, \theta_{1}, \theta_{2}, \theta_{3}\right)$ is an $S$-3-system.

Proof. Obviously $\theta_{i}$ is an equivalence relation, $i=1,2,3$. Also $\theta_{1} \cap \theta_{2}=0$ and $\theta_{1} \cdot \theta_{2}=1$ is clear. If $(x, y)\left(\theta_{1} \cap \theta_{3}\right)\left(x^{\prime}, y^{\prime}\right)$, then $x \cdot y=x^{\prime} \cdot y^{\prime}=x \cdot y^{\prime}$; since $\mathfrak{Q}$ is a quasigroup this implies
$y=y^{\prime}$. Thus, $\theta_{1} \cap \theta_{3}=0$. Let $(x, y)$ and $\left(x^{\prime}, y^{\prime}\right)$ be arbitrary elements of $S$. There exists $y^{\prime \prime} \in S$ such that $x \cdot y^{\prime \prime}=x^{\prime} \cdot y^{\prime}$. Hence, $(x, y) \theta_{1}\left(x, y^{\prime \prime}\right) \theta_{3}\left(x^{\prime}, y^{\prime}\right)$ so that we have $\theta_{1} \cdot \theta_{3}=1$. As the situation with $\theta_{2}$ is analogous this concludes the proof.

For the inverse process start with an S-3-system $\left(S, \theta_{1}, \theta_{2}, \theta_{3}\right)$. Since $\theta_{i} \cdot \theta_{j}=1$ and $\theta_{i} \cap \theta_{j}=0$ for $i \neq j$, we have that $S \cong S / \theta_{i} \times S / \theta_{j}$ as sets (see also Lemma 2.4.5). Hence,

$$
S / \theta_{1} \times S / \theta_{2} \cong S / \theta_{1} \times S / \theta_{3} \cong S / \theta_{2} \times S / \theta_{3}
$$

and thus $S / \theta_{1} \cong S / \theta_{2} \cong S / \theta_{3}$. Therefore, if we set $Q=\left|S / \theta_{1}\right|$, we know there are bijections $f_{i}: S / \theta_{i} \rightarrow Q, i=1,2,3$. Define a function

$$
\begin{array}{lccc}
g: & S & \rightarrow & Q \times Q \\
& s & \mapsto & \left(f_{1}\left([s]_{\theta_{1}}\right), f_{2}\left([s]_{\theta_{2}}\right)\right)
\end{array}
$$

Then $g$ is a bijection: For it is onto since in the corresponding Äquivalenzklassengeometrie two non-parallel lines have an intersection point and it is one-one since this intersection is unique. Fix an arbitrary element $e \in S$ and an arbitrary element $1 \in Q$. We may suppose we have chosen $f_{1}, f_{2}$ such that $f_{1}\left([e]_{\theta_{1}}\right)=f_{2}\left([e]_{\theta_{2}}\right)=1$. Furthermore, we suppose that if $a \theta_{1} e, b \theta_{2} e$ and $(a, b) \notin \theta_{3}$ then $f_{2}\left([a]_{\theta_{2}}\right) \neq f_{1}\left([b]_{\theta_{1}}\right)$. This is legitimate since there are $|Q|^{2}-|Q|$ possibilities to choose an ordered pair of two distinct equivalence classes in $\theta_{3}$; but the other assumptions on $a$ and $b$ already uniquely determine the representatives of those equivalence classes. Hence, there are $|Q|^{2}-|Q|$ possibilities to choose $a$ and $b$ satisfying all conditions which is exactly the number of ordered pairs of unequal values of $f_{1}$ and $f_{2}$. Note that the assumption implies $g^{-1}(1, x) \theta_{3} g^{-1}(x, 1)$ for all $x \in Q$. Define a binary operation • on $Q$ in the following way: For $x, y$ in $Q$ set $s=g^{-1}(x, y)$. Let $t$ be the intersection of the $\theta_{3}$-line through $s$ with the $\theta_{2}$-line through $e$; then $z=x \cdot y=g(t)$. More formally,

$$
x \cdot y=z \leftrightarrow g^{-1}(x, y) \theta_{3} g^{-1}(z, 1)
$$

2.3.11 Definition. A loop is a quasigroup $\mathfrak{L}=(L, \cdot)$ which has an element $1 \in L$ such that $x \cdot 1=1 \cdot x=x$ for all $x \in L$.
2.3.12 Lemma. $\mathfrak{Q}=(Q, \cdot, 1)$ is a loop.

Proof. $x \cdot 1=x$ since trivially $g^{-1}(x, 1) \theta_{3} g^{-1}(x, 1) ; 1 \cdot x=x$ since we chose $g$ such that $g^{-1}(1, x) \theta_{3} g^{-1}(x, 1)$. To find the right-side inverse of an element $x \in Q$, let $s$ be the intersection of the $\theta_{1}$-line $f_{1}^{-1}(x)$ with the $\theta_{3}$-line through $e$. Then for $y=f_{2}\left([s]_{\theta_{2}}\right)$ we have $x \cdot y=1$ : $g^{-1}(x, y)=s \theta_{3} e=g^{-1}(1,1)$. The left-side inverse can be found in a similar way.

In the following, we will identify $S$ with $Q \times Q$. Then one can easily verify

$$
\begin{array}{ccc}
(x, y) \theta_{1}\left(x^{\prime}, y^{\prime}\right) & \leftrightarrow & x=x^{\prime} \\
(x, y) \theta_{2}\left(x^{\prime}, y^{\prime}\right) & \leftrightarrow & y=y^{\prime} \\
(x, y) \theta_{3}\left(x^{\prime}, y^{\prime}\right) & \leftrightarrow & x \cdot y=x^{\prime} \cdot y^{\prime}
\end{array}
$$

Hence, if we start out with an S-3-system, construct a loop as shown before, and construct from that an S-3-system again, we end up with the system we started with. We summarize this connection in the following theorem.
2.3.13 Theorem. Let $\mathcal{S}=\left(S, \theta_{1}, \theta_{2}, \theta_{3}\right)$ be an $S$-3-system and let $e \in S$ arbitrary. Then there exist a loop $\mathfrak{L}=(L, \cdot, 1)$ and a bijection $g: L \times L \rightarrow S$ such that $e=g(1,1)$ and for all $x, y, x^{\prime}, y^{\prime} \in L$ we have

$$
\begin{array}{rcc}
(x, y) \theta_{1}\left(x^{\prime}, y^{\prime}\right) & \leftrightarrow & x=x^{\prime} \\
(x, y) \theta_{2}\left(x^{\prime}, y^{\prime}\right) & \leftrightarrow & y=y^{\prime} \\
(x, y) \theta_{3}\left(x^{\prime}, y^{\prime}\right) & \leftrightarrow & x \cdot y=x^{\prime} \cdot y^{\prime}
\end{array}
$$

if we identify the elements of $S$ with those of $L \times L$ via $g$.
2.3.14 Remark. Note that $\theta_{3}$ need not be a congruence of $\mathfrak{L} \times \mathfrak{L}$ whereas $\theta_{1}$ and $\theta_{2}$ obviously are.

Now let us return to our algebra $\mathfrak{A}$. Since $\mathfrak{A}$ generates a congruence permutable variety, there exists a Mal'cev term on $\mathfrak{A}$, that is, there exists a ternary term $p$ satisfying the equations $p(x, x, y)=y$ and $p(x, y, y)=x$. In our case, $p$ is unique:
2.3.15 Theorem. Let $\mathcal{S}=\left(S, \theta_{1}, \theta_{2}, \theta_{3}\right)$ be an $S$-3-system and let $p$ be a Mal'cev function on $S$ preserving $\theta_{1}, \theta_{2}$, and $\theta_{3}$. Then $p$ is uniquely determined.

Proof. Let $x, y, z \in S$ be given. If $x=y$ or $y=z$ then $p(x, y, z)$ is determined by the equations of a Mal'cev function. Suppose $x \neq y$ and $y \neq z$. Assume first that $x$ and $y$ lie on one line $l_{1}$ and $y$ and $z$ lie on a line $l_{2}$ and $l_{1} \neq l_{2}$. Since $y$ is the intersection of those lines, for some $i \neq k$ we have $l_{1}=[y]_{\theta_{i}}$ and $l_{2}=[y]_{\theta_{k}}$ so that by compatibility $p(x, y, z) \theta_{i} p(x, x, z)=z$ and $p(x, y, z) \theta_{k} p(x, y, y)=x$. Hence, $p(x, y, z)$ is the intersection of the $\theta_{i}$-line through $z$ with the $\theta_{k}$-line through $x$ so that it must be unique. In a next step, assume that $x, y, z$ lie on one line $l$ and say without loss of generality $l$ is a $\theta_{1}$-line. Denote by $x^{\prime}$ the intersection of the $\theta_{2}$-line through $y$ with the $\theta_{3}$-line through $x$. As $x^{\prime}$ and $y$ lie on one line and $y$ and $z$ on another line, we know from the first step of the proof that $p\left(x^{\prime}, y, z\right)$ is uniquely determined. Since $x, y, z$ lie on one $\theta_{1}$-line we have $x=p(x, x, x) \theta_{1} p(x, y, z)$; hence, $p(x, y, z)$ lies on $l$ as well. But $x \theta_{3} x^{\prime}$ implies $p(x, y, z) \theta_{3} p\left(x^{\prime}, y, z\right)$. Thus, $p(x, y, z)$ is the intersection of the $\theta_{3}$-line through $p\left(x^{\prime}, y, z\right)$ with $l$ and so it is uniquely determined. To finish the proof, let $x, y, z$ be arbitrary. Consider an arbitrary $\theta_{1}$-line $l_{1}$ and an arbitrary $\theta_{2}$-line $l_{2}$. Denote the intersections of the $\theta_{2}$-lines through $x, y, z$ with $l_{1}$ by $x^{\prime}, y^{\prime}, z^{\prime}$ and the intersections of the $\theta_{1}$-lines through $x, y, z$ with $l_{2}$ by $x^{\prime \prime}, y^{\prime \prime}, z^{\prime \prime}$. By the second step of our proof, $p\left(x^{\prime}, y^{\prime}, z^{\prime}\right)$ and $p\left(x^{\prime \prime}, y^{\prime \prime}, z^{\prime \prime}\right)$ are uniquely determined. Clearly, $p(x, y, z) \theta_{2} p\left(x^{\prime}, y^{\prime}, z^{\prime}\right)$ and $p(x, y, z) \theta_{1} p\left(x^{\prime \prime}, y^{\prime \prime}, z^{\prime \prime}\right)$. Hence, $p(x, y, z)$ is the unique intersection of the $\theta_{2}$-line through $p\left(x^{\prime}, y^{\prime}, z^{\prime}\right)$ with the $\theta_{1}$-line through $p\left(x^{\prime \prime}, y^{\prime \prime}, z^{\prime \prime}\right)$.
2.3.16 Corollary. Let $\mathcal{S}=\left(S, \theta_{1}, \theta_{2}, \theta_{3}\right)$ be an $S$-3-system. If $p$ is a Mal'cev function on $S$ compatible with $\theta_{1}, \theta_{2}, \theta_{3}$, then it satisfies the equation $p(x, y, z)=p(z, y, x)$.

Proof. Set $\tilde{p}(x, y, z)=p(z, y, x)$. Then $\tilde{p}$ is apparently a Mal'cev function on $S$ preserving $\theta_{1}, \theta_{2}, \theta_{3}$. Therefore, it must equal $p$ and so $\tilde{p}(x, y, z)=p(z, y, x)=p(x, y, z)$ for all $x, y, z \in$ $S$.
2.3.17 Lemma. Let $\mathcal{S}=\left(S, \theta_{1}, \theta_{2}, \theta_{3}\right)$ be an $S$-3-system. If there is a compatible Mal'cev function on $S$, then the loop associated with $\mathcal{S}$ satisfies:

$$
\left(x_{1} \cdot y_{1}=x_{2} \cdot y_{2} \wedge x_{1} \cdot y_{3}=x_{2} \cdot y_{4} \wedge x_{3} \cdot y_{1}=x_{4} \cdot y_{2}\right) \rightarrow x_{3} \cdot y_{3}=x_{4} \cdot y_{4}
$$

Proof. Recall that in terms of the S-3-system and its congruence $\theta_{3}$ our hypothesis says

$$
\left(x_{1}, y_{1}\right) \theta_{3}\left(x_{2}, y_{2}\right) \wedge\left(x_{1}, y_{3}\right) \theta_{3}\left(x_{2}, y_{4}\right) \wedge\left(x_{3}, y_{1}\right) \theta_{3}\left(x_{4}, y_{2}\right)
$$

Since $p$ is compatible with $\theta_{1}$ and $\theta_{2}$, it satisfies the Mal'cev conditions componentwise. Hence, $\left(x_{3}, y_{3}\right)=p\left(\left(x_{1}, y_{3}\right),\left(x_{1}, y_{1}\right),\left(x_{3}, y_{1}\right)\right) \theta_{3} p\left(\left(x_{2}, y_{4}\right),\left(x_{2}, y_{2}\right),\left(x_{4}, y_{2}\right)\right)=\left(x_{4}, y_{4}\right)$.
2.3.18 Lemma. Let $\mathcal{S}$ be an $S$-3-system with a compatible Mal'cev function. Then the loop $\mathfrak{L}$ associated with $\mathcal{S}$ is associative, i.e. a group.

Proof. The previous lemma applies; so to check the associative law for arbitrary $x, y, z \in L$ set $x_{1}=y_{2}=1, x_{2}=y, x_{3}=x, x_{4}=x \cdot y, y_{1}=y, y_{3}=y \cdot z, y_{4}=z$. Then all hypotheses of the lemma are satisfied and it yields $x \cdot(y \cdot z)=(x \cdot y) \cdot z$.

So we know that if we have a Mal'cev operation compatible with an S-3-system, the associated loop is in fact a group. We will show now that this group is even abelian.
2.3.19 Lemma. Let the $S$-3-system $\mathcal{S}$ admit the Mal'cev function $p$. Then we can calculate $p$ by

$$
p\left(\left(x_{1}, y_{1}\right),\left(x_{2}, y_{2}\right),\left(x_{3}, y_{3}\right)\right)=\left(x_{1} \cdot x_{2}^{-1} \cdot x_{3}, y_{1} \cdot y_{2}^{-1} \cdot y_{3}\right)
$$

Proof. We will calculate $p(x, y, z)$ following the construction of that point in the proof of Theorem 2.3.15. Set $x=\left(x_{1}, y_{1}\right), y=\left(x_{2}, y_{2}\right)$ and $z=\left(x_{3}, y_{3}\right)$. Let $l_{1}$ be the $\theta_{1}$-line and $l_{2}$ be the $\theta_{2}$-line through $x$. Then, using the same notation as in that proof, we have $x^{\prime}=\left(x_{1}, y_{1}\right)$, $y^{\prime}=\left(x_{1}, y_{2}\right), z^{\prime}=\left(x_{1}, y_{3}\right)$ and $x^{\prime \prime}=\left(x_{1}, y_{1}\right), y^{\prime \prime}=\left(x_{2}, y_{1}\right), z^{\prime \prime}=\left(x_{3}, y_{1}\right)$. Thus, if we write $p$ also for the functions $p$ induces on the components,

$$
\begin{aligned}
p\left(x^{\prime}, y^{\prime}, z^{\prime}\right) & =p\left(\left(x_{1}, y_{1}\right),\left(x_{1}, y_{2}\right),\left(x_{1}, y_{3}\right)\right) \\
& =\left(p\left(x_{1}, x_{1}, x_{1}\right), p\left(y_{1}, y_{2}, y_{3}\right)\right) \\
& =:\left(x_{1}, \bar{p}\right)
\end{aligned}
$$

and similarly $p\left(x^{\prime \prime}, y^{\prime \prime}, z^{\prime \prime}\right)=\left(p\left(x_{1}, x_{2}, x_{3}\right), y_{1}\right)=:\left(\overline{\bar{p}}, y_{1}\right)$. Since $p(x, y, z) \theta_{2} p\left(x^{\prime}, y^{\prime}, z^{\prime}\right)$ and $p(x, y, z) \theta_{1} p\left(x^{\prime \prime}, y^{\prime \prime}, z^{\prime \prime}\right)$ we get $p(x, y, z)=(\overline{\bar{p}}, \bar{p})$. For the computation of $p\left(x^{\prime}, y^{\prime}, z^{\prime}\right)$ we can use the second step in the proof of 2.3.15 since $x^{\prime}, y^{\prime}, z^{\prime}$ lie on one line $l_{1}$ : Let $s=\left(u, y_{2}\right)$ be the intersection of the $\theta_{3}$-line through $x^{\prime}$ with the $\theta_{2}$-line through $y^{\prime}$; the definition of $\theta_{3}$
immediately yields the equation $x_{1} \cdot y_{1}=u \cdot y_{2}$. Let $t=\left(u, y_{3}\right)$ be the intersection of the $\theta_{1}$-line through $s$ with the $\theta_{2}$-line through $z^{\prime}$. Then, as $p\left(x^{\prime}, y^{\prime}, z^{\prime}\right)=\left(x_{1}, \bar{p}\right)$ is the intersection of the $\theta_{3}$-line through $t$ with $l_{1}$, we get $x_{1} \cdot \bar{p}=u \cdot y_{3}$ which we can solve to

$$
\bar{p}=x_{1}^{-1} \cdot u \cdot y_{3}=x_{1}^{-1} \cdot x_{1} \cdot y_{1} \cdot y_{2}^{-1} \cdot y_{3}=y_{1} \cdot y_{2}^{-1} \cdot y_{3} .
$$

Similarly, $\overline{\bar{p}}=x_{1} \cdot x_{2}^{-1} \cdot x_{3}$ and so the proof is complete.
2.3.20 Corollary. If $\mathcal{S}=\left(S, \theta_{1}, \theta_{2}, \theta_{3}\right)$ is an $S$-3-system which allows a Mal'cev function on $S$, then the associated group $\mathfrak{G}$ is abelian.

Proof. Combining Lemma 2.3 .19 with Corollary 2.3 .16 yields $\left(x_{1} \cdot x_{2}^{-1} \cdot x_{3}, y_{1} \cdot y_{2}^{-1} \cdot y_{3}\right)=$ $\left(x_{3} \cdot x_{2}^{-1} \cdot x_{1}, y_{1} \cdot y_{2}^{-1} \cdot y_{3}\right)$ for all $x_{i}, y_{i} \in G, i=1,2,3$. Therefore, if we set $x_{2}=1$, we have $x_{1} \cdot x_{3}=x_{3} \cdot x_{1}$.

Since we know now that we are dealing an abelian group, we will change our notation to an additive one; that is, we will write + for the binary group operation and 0 for the neutral element. Furthermore, we will identify the base set $S$ of an S-3-system that admits a Mal'cev operation with $G \times G$, where $\mathfrak{G}=(G,+, 0,-)$ is the associated abelian group. Observe that for the Mal'cev operation $p$ on $G \times G$ we have $p(x, y, z)=x-y+z$, where + is calculated componentwise.

Now let $f:(G \times G)^{n} \rightarrow G$ be an $n$-ary function on $G \times G$ compatible with $\theta_{1}, \theta_{2}, \theta_{3}$. Since $f$ is compatible with $\theta_{1}$ and $\theta_{2}$, it is the product of two mappings $f_{1}, f_{2}: G^{n} \rightarrow G$, i.e.

$$
f\left(\left(x_{1}, y_{1}\right), \ldots,\left(x_{n}, y_{n}\right)\right)=\left(f_{1}\left(x_{1}, \ldots, x_{n}\right), f_{2}\left(y_{1}, \ldots, y_{n}\right)\right)
$$

2.3.21 Lemma. For $x, y, x^{\prime}, y^{\prime} \in G^{n}$ the following holds:
(i) $x+y=x^{\prime}+y^{\prime} \rightarrow f_{1}(x)+f_{2}(y)=f_{1}\left(x^{\prime}\right)+f_{2}\left(y^{\prime}\right)$
(ii) $f_{1}(x)+f_{2}(0)=f_{1}(0)+f_{2}(x)$
(iii) $f_{1}(x)+f_{2}(y)=f_{1}(x+y)+f_{2}(0)$

Proof. The hypothesis of (i) says that for $1 \leq k \leq n, x_{k}+y_{k}=x_{k}^{\prime}+y_{k}^{\prime}$. Thus, by definition of $\theta_{3},\left(x_{k}, y_{k}\right) \theta_{3}\left(x_{k}^{\prime}, y_{k}^{\prime}\right)$ for $1 \leq k \leq n$ so that by the compatibility of $f$ with $\theta_{3}$, $f\left(\left(x_{1}, y_{1}\right), \ldots,\left(x_{n}, y_{n}\right)\right) \theta_{3} f\left(\left(x_{1}^{\prime}, y_{1}^{\prime}\right), \ldots,\left(x_{n}^{\prime}, y_{n}^{\prime}\right)\right)$. Hence, $\left(f_{1}(x), f_{2}(y)\right) \theta_{3}\left(f_{1}\left(x^{\prime}\right), f_{2}\left(y^{\prime}\right)\right)$ and so $f_{1}(x)+f_{2}(y)=f_{1}\left(x^{\prime}\right)+f_{2}\left(y^{\prime}\right)$. (ii) is trivial with (i). For (iii) set $x^{\prime}=x+y$ and $y^{\prime}=0$ and apply (i).
2.3.22 Definition. Let $A$ be a set and $f: A^{n} \rightarrow A$ be an $n$-ary function on $A$. If a binary operation + can be defined on $A$ such that $(A,+)$ is an abelian group and for all $x, y \in A^{n}$ we have

$$
f(x)+f(y)=f(x+y)+f(0)
$$

then we say that $f$ is affine with respect to $(A,+)$. An algebra $\mathfrak{A}$ is affine iff every fundamental operation is affine with respect to the same abelian group over $A$.
2.3.23 Remark. It is obvious that an algebra $\mathfrak{A}$ is affine if and only if there is an affine relation $\rho \subseteq A^{4}$ as defined in the first chapter which is preserved by all operations of $\mathfrak{A}$.
2.3.24 Lemma. Let $\mathcal{S}=\left(S, \theta_{1}, \theta_{2}, \theta_{3}\right)$ be an $S$-3-system with a compatible Mal'cev function. Then every mapping on $S$ which is compatible with $\theta_{1}, \theta_{2}, \theta_{3}$ is affine with respect to $\mathfrak{G} \times \mathfrak{G}$, where $\mathfrak{G}$ is the abelian group associated with $\mathcal{S}$.

Proof. Suppose $f: S^{n} \rightarrow S$ is compatible with $\theta_{1}, \theta_{2}, \theta_{3}$. Then for $x, y \in S$, if we write $x=\left(x^{\prime}, x^{\prime \prime}\right), y=\left(y^{\prime}, y^{\prime \prime}\right)$, we can compute by the previous lemma:

$$
\begin{aligned}
f(x)+f(y) & =\left(f_{1}\left(x^{\prime}\right), f_{2}\left(x^{\prime \prime}\right)\right)+\left(f_{1}\left(y^{\prime}\right), f_{2}\left(y^{\prime \prime}\right)\right) \\
& =\left(f_{1}\left(x^{\prime}\right)+f_{1}\left(y^{\prime}\right), f_{2}\left(x^{\prime \prime}\right)+f_{2}\left(y^{\prime \prime}\right)\right) \\
& =\left(f_{1}\left(x^{\prime}\right)+f_{1}(0)+f_{2}\left(y^{\prime}\right)-f_{2}(0), f_{2}\left(x^{\prime \prime}\right)+f_{2}(0)+f_{1}\left(y^{\prime \prime}\right)-f_{1}(0)\right) \\
& =\left(f_{1}\left(x^{\prime}+y^{\prime}\right)+f_{1}(0), f_{2}\left(x^{\prime \prime}+y^{\prime \prime}\right)+f_{2}(0)\right) \\
& =\left(f_{1}\left(x^{\prime}+y^{\prime}\right), f_{2}\left(x^{\prime \prime}+y^{\prime \prime}\right)\right)+\left(f_{1}(0), f_{2}(0)\right) \\
& =f(x+y)+f(0) .
\end{aligned}
$$

In terms of algebras we have established:
2.3.25 Theorem. Let $\mathfrak{A}$ be an algebra in a congruence permutable variety and let $p$ be a Mal'cev term of $\mathfrak{A}$. If $\mathcal{M}_{3}$ is a 0-1-sublattice of Con $(\mathfrak{A})$, then there is an abelian group $\mathfrak{G}=(G,+, 0)$ such that $\mathfrak{A}$ is isomorphic as sets to $\mathfrak{G} \times \mathfrak{G}$ and such that the following holds: $p(x, y, z)=x-y+z$ and every term operation $f$ is affine with respect to $\mathfrak{G} \times \mathfrak{G}$ and of the form $f_{1} \times f_{2}$ where $f_{1}, f_{2}: G^{n} \rightarrow G$ if $f$ is $n$-ary.

The connection to skew congruences is the following:
2.3.26 Theorem. Let $\mathfrak{A}$ be a simple algebra in a congruence permutable variety. If $\mathfrak{A} \times \mathfrak{A}$ has a skew congruence, then $\mathfrak{A}$ is affine.

Proof. If $\mathfrak{A} \times \mathfrak{A}$ has a skew congruence $\theta$, then $\theta$ is by our discussion at the beginning of this section a complement of $\operatorname{ker}\left(\pi_{1}^{2}\right)$ and of $\operatorname{ker}\left(\pi_{2}^{2}\right)$. Hence, the last theorem applies to $\mathfrak{A} \times \mathfrak{A}$. But it follows from the construction of $\mathfrak{G} \times \mathfrak{G}$ by means of the congruences $\operatorname{ker}\left(\pi_{1}^{2}\right)$ and $\operatorname{ker}\left(\pi_{2}^{2}\right)$ that the canonical coordinate representation of an element of $\mathfrak{A} \times \mathfrak{A}$ is exactly the same as the representation with respect to the factorization $\mathfrak{G} \times \mathfrak{G}$. We conclude that every term operation on $\mathfrak{A}$ is affine with respect to $\mathfrak{G}$ as the corresponding term operation on $\mathfrak{A} \times \mathfrak{A}$ is affine with respect to $\mathfrak{G} \times \mathfrak{G}$. Thus, $\mathfrak{A}$ is affine with respect to $\mathfrak{G}$.

Our next goal is to show that for a simple non-trivial affine algebra the underlying abelian group is in fact a $p$-group for some prime $p$. Define for arbitrary $n \geq 1$ a binary relation $\Delta_{n}$ on $\mathfrak{A}$ by

$$
x \Delta_{n} y \leftrightarrow n(x-y)=0 .
$$

2.3.27 Lemma. Let $\mathfrak{A}$ be an affine algebra. Then $\Delta_{n}$ is a congruence on $\mathfrak{A}$ for all $n \geq 1$.

Proof. For an arbitrary $k$-ary fundamental operation of $\mathfrak{A}$, if $x_{i} \Delta_{n} y_{i}$ for $1 \leq i \leq k$, we have

$$
\begin{aligned}
n\left(f\left(x_{1}, \ldots, x_{k}\right)-f\left(y_{1}, \ldots, y_{k}\right)\right) & =n\left(f\left(x_{1}-y_{1}, \ldots, x_{k}-y_{k}\right)-f(0, \ldots, 0)\right) \\
& =n f\left(x_{1}-y_{1}, \ldots, x_{k}-y_{k}\right)-n f(0, \ldots, 0) \\
& =f\left(n\left(x_{1}-y_{1}\right), \ldots, n\left(x_{k}-y_{k}\right)\right)-f(0, \ldots, 0) \\
& =f(0, \ldots, 0)-f(0, \ldots, 0) \\
& =0
\end{aligned}
$$

2.3.28 Lemma. Let $\mathfrak{A}$ be a simple non-trivial affine algebra. Then the underlying group $\mathfrak{G}$ is either torsion-free or a p-group for some prime $p$.

Proof. Suppose $\mathfrak{G}$ is not torsion-free. Then there exists a smallest positive number $p$ such that $p a=0$ for some $a \in G, a \neq 0$. Obviously $p$ is a prime. Consider the congruence $\Delta_{p}$ of the previous lemma. Since $a \Delta_{p} 0$ and $a \neq 0, \Delta_{p}$ must equal $1 \in \operatorname{Con}(\mathfrak{A})$ as $\mathfrak{A}$ is simple. But then it readily follows that $p a=0$ for all $a \in G$.
2.3.29 Theorem. Let $\mathfrak{A}$ be a simple non-trivial algebra in a permutable variety. If $\mathfrak{A} \times \mathfrak{A}$ has a skew congruence, then $\mathfrak{A}$ is affine with respect to a torsion-free abelian group or with respect to an abelian p-group.

Proof. Follows from Theorem 2.3.26 and the previous lemma.
The following theorem (see H. Werner [13] for a slightly more general result) tells us that if we have a skew congruence on a higher power of $\mathfrak{A}$ we can still use our results and $\mathfrak{A}$ must be prime affine as well.
2.3.30 Theorem. Let $\mathfrak{A}_{1}, \ldots, \mathfrak{A}_{n}$ be algebras in a permutable variety. Then $\mathfrak{A}_{1} \times \ldots \times \mathfrak{A}_{n}$ has a skew congruence if and only if $\mathfrak{A}_{i} \times \mathfrak{A}_{j}$ has a skew congruence for some $1 \leq i, j \leq n$ with $i \neq j$.

Before we can prove the theorem, we need a definition. For $\theta \in \operatorname{Con}(\mathfrak{A} \times \mathfrak{B})$ and $a \in A$, we define an equivalence relation $\theta_{a}$ on $\mathfrak{B}$ by

$$
\theta_{a}=\left\{\left(b_{1}, b_{2}\right) \mid\left(a, b_{1}\right) \theta\left(a, b_{2}\right)\right\}
$$

In a congruence permutable variety, $\theta_{a}$ is a congruence relation and independent of the choice of $a$ :
2.3.31 Lemma. Let $\mathfrak{A}, \mathfrak{B}$ be algebras in a variety with permutable congruences and let $\theta \in$ $\operatorname{Con}(\mathfrak{A} \times \mathfrak{B})$. Then for all $a, a^{\prime} \in A$ and all $b \in B$
(i) $\theta_{a}=\theta_{a^{\prime}}$
(ii) $\theta_{a} \in \operatorname{Con}(\mathfrak{B})$
(iii) $\theta_{b} \times \theta_{a} \subseteq \theta$

Proof. To prove (i), denote the Mal'cev term of the variety by $p$. We have to show for arbitrary $a, a^{\prime} \in A$ and $b, b^{\prime} \in B$ that $(a, b) \theta\left(a, b^{\prime}\right)$ implies $\left(a^{\prime}, b\right) \theta\left(a^{\prime}, b^{\prime}\right)$. But as trivially $\left(a^{\prime}, b^{\prime}\right) \theta\left(a^{\prime}, b^{\prime}\right)$ and $(a, b) \theta(a, b)$, application of $p$ to the three elements of $\theta$ yields

$$
\left(a^{\prime}, b^{\prime}\right)=\left(p\left(a^{\prime}, a, a\right), p\left(b^{\prime}, b, b\right)\right) \theta\left(p\left(a^{\prime}, a, a\right), p\left(b^{\prime}, b^{\prime}, b\right)\right)=\left(a^{\prime}, b\right)
$$

so that (i) is indeed true. Now (ii) is an immediate consequence of (i) since if $f$ is an $n$-ary operation of $\mathfrak{B}$ and $\left(b_{i}, b_{i}^{\prime}\right) \in \theta_{a}, 1 \leq i \leq n$, then $\left(f\left(b_{1}, \ldots, b_{n}\right), f\left(b_{1}^{\prime}, \ldots, b_{n}^{\prime}\right)\right) \in \theta_{f(a, \ldots, a)}=$ $\theta_{a}$. For (iii), let $\left((c, d),\left(c^{\prime}, d^{\prime}\right)\right) \in \theta_{b} \times \theta_{a}$ be given. By (i), we can replace $b$ by $d$ and $a$ by $c^{\prime}$. Then we get by the definitions of $\theta_{d}$ and $\theta_{c^{\prime}}$ that $(c, d) \theta\left(c^{\prime}, d\right)$ and $\left(c^{\prime}, d\right) \theta\left(c^{\prime}, d^{\prime}\right)$ so that $\left((c, d),\left(c^{\prime}, d^{\prime}\right)\right) \in \theta$.
2.3.32 Corollary. Let $\mathfrak{A}, \mathfrak{B}$ be algebras in a variety with permutable congruences and let $\theta \in$ $\operatorname{Con}(\mathfrak{A} \times \mathfrak{B})$. Then the following conditions are equivalent:
(i) $\theta$ is not skew.
(ii) $\theta_{b} \times \theta_{a}=\theta$.
(iii) $(a, b) \theta\left(a^{\prime}, b^{\prime}\right)$ implies $(a, b) \theta\left(a, b^{\prime}\right)$.
$\left(i i i^{*}\right)(a, b) \theta\left(a^{\prime}, b^{\prime}\right)$ implies $(a, b) \theta\left(a^{\prime}, b\right)$.
Proof. (i) $\Rightarrow$ (ii): Let $\theta=\theta_{1} \times \theta_{2}$. If $\left(a, a^{\prime}\right) \in \theta_{1}$, then since trivially $(b, b) \in \theta_{2}$ we have $\left((a, b),\left(a^{\prime}, b\right)\right) \in \theta_{1} \times \theta_{2}=\theta$ so that $\left(a, a^{\prime}\right) \in \theta_{b}$. As the situation for $\theta_{2}$ is the same we conclude $\theta_{b} \times \theta_{a} \supseteq \theta$ and so $\theta_{b} \times \theta_{a}=\theta$. (ii) $\Rightarrow$ (i) and (ii) $\Rightarrow$ (iii) are trivial. To see that (iii) implies (ii), let $(a, b) \theta\left(a^{\prime}, b^{\prime}\right)$. By hypothesis, $(a, b) \theta\left(a, b^{\prime}\right)$ so that $\left(b, b^{\prime}\right) \in \theta_{a}$. But the transitivity of $\theta$ implies $\left(a^{\prime}, b^{\prime}\right) \theta\left(a, b^{\prime}\right)$. Hence, $\left(a, a^{\prime}\right) \in \theta_{b^{\prime}}=\theta_{b}$.

Now we are ready to prove the theorem.
Proof of Theorem 2.3.30. We have to show that if $\mathfrak{A}_{1} \times \ldots \times \mathfrak{A}_{n}$ has a skew congruence then $\mathfrak{A}_{i} \times \mathfrak{A}_{j}$ has a skew congruence for some $1 \leq i, j \leq n$ with $i \neq j$ as the other direction is obvious. To achieve this we will prove for algebras $\mathfrak{A}, \mathfrak{B}, \mathfrak{C}$ of our variety that if $\mathfrak{A} \times \mathfrak{B}, \mathfrak{A} \times \mathfrak{C}$, and $\mathfrak{B} \times \mathfrak{C}$ have only factor congruences, then $\mathfrak{A} \times \mathfrak{B} \times \mathfrak{C}$ can have only factor congruences as well. The
rest will follow by induction. Let $\theta \in \operatorname{Con}(\mathfrak{A} \times \mathfrak{B} \times \mathfrak{C})$. Define $\phi_{A}=\theta_{(b, c)}, \phi_{B}=\theta_{(a, c)}$ and $\phi_{C}=\theta_{(a, b)}$. By Lemma 2.3.31 (iii), $\phi_{A} \times \phi_{B} \times \phi_{C} \subseteq \theta$.

Claim. $(a, b, c) \theta\left(a^{\prime}, b^{\prime}, c^{\prime}\right) \wedge a \phi_{A} a^{\prime} \rightarrow b \phi_{B} b^{\prime} \wedge c \phi_{C} c^{\prime}$.
Proof. Since $\phi_{A} \times \phi_{B} \times \phi_{C} \subseteq \theta,(a, b, c) \theta\left(a^{\prime}, b^{\prime}, c^{\prime}\right) \theta\left(a, b^{\prime}, c^{\prime}\right)$. Thus, $(b, c) \theta_{a}\left(b^{\prime}, c^{\prime}\right)$. As $\theta_{a}$ is not skew, $(b, c) \theta_{a}\left(b, c^{\prime}\right)$ by Corollary 2.3.32 (iii). Therefore, $(a, b, c) \theta\left(a, b, c^{\prime}\right)$ and so $\left(c, c^{\prime}\right) \in \phi_{C}$. Hence, $(a, b, c) \theta\left(a^{\prime}, b^{\prime}, c^{\prime}\right) \theta\left(a, b^{\prime}, c\right)$ so that $\left(b, b^{\prime}\right) \in \phi_{B}$.

Denote by $\pi$ the projection of $A \times B \times C$ onto $A \times B$. Then $\tilde{\theta}=\pi[\theta]$ is a congruence on $\mathfrak{A} \times \mathfrak{B}$ : The only property which is not obvious is the transitivity of $\tilde{\theta}$. Assume $(a, b) \tilde{\theta}\left(a^{\prime}, b^{\prime}\right)$ and $\left(a^{\prime}, b^{\prime}\right) \tilde{\theta}\left(a^{\prime \prime}, b^{\prime \prime}\right)$; then there exist $c, c^{\prime}, d, d^{\prime} \in C$ such that $(a, b, c) \theta\left(a^{\prime}, b^{\prime}, c^{\prime}\right)$ and $\left(a^{\prime}, b^{\prime}, d\right) \theta\left(a^{\prime \prime}, b^{\prime \prime}, d^{\prime}\right)$. If $p$ is a Mal'cev term of the variety, then since trivially $\left(a^{\prime \prime}, b^{\prime \prime}, d^{\prime}\right) \theta\left(a^{\prime \prime}, b^{\prime \prime}, d^{\prime}\right)$ we get

$$
\left(p\left(a, a^{\prime \prime}, a^{\prime \prime}\right), p\left(b, b^{\prime \prime}, b^{\prime \prime}\right), p\left(c, d^{\prime}, d^{\prime}\right)\right) \theta\left(p\left(a^{\prime}, a^{\prime}, a^{\prime \prime}\right), p\left(b^{\prime}, b^{\prime}, b^{\prime \prime}\right), p\left(c^{\prime}, d, d^{\prime}\right)\right)
$$

This yields $(a, b, c) \theta\left(a^{\prime \prime}, b^{\prime \prime}, p\left(c^{\prime}, d, d^{\prime}\right)\right)$ and so $(a, b) \tilde{\theta}\left(a^{\prime \prime}, b^{\prime \prime}\right)$. Hence, $\tilde{\theta}$ is a congruence which is not skew by assumption. Now if $(a, b, c) \theta\left(a^{\prime}, b^{\prime}, c^{\prime}\right)$, then $(a, b) \tilde{\theta}\left(a^{\prime}, b^{\prime}\right)$ so that also $(a, b) \tilde{\theta}\left(a, b^{\prime}\right)$. Hence, there exist $d, d^{\prime} \in C$ such that $(a, b, d) \theta\left(a, b^{\prime}, d^{\prime}\right)$. Since trivially $(a, a) \in \phi_{A}$, by our claim we get that $\left(b, b^{\prime}\right) \in \phi_{B}$ and again by the claim that $\left(a, a^{\prime}\right) \in \phi_{A}$ and $\left(c, c^{\prime}\right) \in \phi_{C}$. Therefore, $\theta \subseteq \phi_{A} \times \phi_{B} \times \phi_{C}$ which is exactly what we wanted to show.
2.3.33 Theorem. Let $\mathfrak{A}$ be a finite simple algebra in a permutable variety. Then $\mathfrak{A}$ is either prime affine or its powers have only (trivial) factor congruences.

Proof. If $\mathfrak{A} \times \mathfrak{A}$ has a skew congruence, then by Theorem 2.3.29 and the fact that $\mathfrak{A}$ is finite we get that $\mathfrak{A}$ is prime affine. Otherwise Theorem 2.3.30 applies and all powers of $\mathfrak{A}$ have only factor congruences.

To use this result in the proof of primality of our algebra $\mathfrak{A}$, recall that since $\{\mathfrak{A}\}$ is a direct factor set, $\mathcal{V}(\mathfrak{A})=H P(\mathfrak{A})$. Now by the last theorem, if a power of $\mathfrak{A}$ has a skew congruence, then $\mathfrak{A}$ is prime affine which is forbidden by Rosenberg's list. Thus, all powers of $\mathfrak{A}$ have only trivial factor congruences: Since $\mathfrak{A}$ is simple, they are products of 0 and $1 \in \operatorname{Con}(\mathfrak{A})$. But this implies also that up to isomorphism $\mathcal{V}(\mathfrak{A})=P(\mathfrak{A})$ and so the variety generated by $\mathfrak{A}$ is obviously congruence distributive.

## $2.4 \mathfrak{A}$ is primal

We will use a special case of a result on semi-primal algebras by A. Foster and A. Pixley in [6] to show that our hypotheses on $\mathfrak{A}$ imply it is primal.
2.4.1 Theorem. Let $\mathfrak{A}$ be an algebra, $1<|\mathfrak{A}|<\aleph_{0}$. Assume also that $\mathfrak{A}$ is simple, has no proper subalgebras and no proper automorphisms and that it generates a congruence permutable and congruence distributive variety. Then $\mathfrak{A}$ is primal.
2.4.2 Remark. Non-trivial congruences are obviously exactly class three in $R B L$; proper subalgebras are central relations and thus in class five of $R B L$. Moreover, assume $\phi$ is a proper automorphism of $\mathfrak{A}$. Now either all cycles of $\phi$ have the same length $n$; then for any prime factor $p$ of $n, \phi^{\frac{n}{p}}$ has only cycles of the same prime length $p$ and hence its graph belongs to class two of $R B L$. If there are cycles of different length, then denote the length of the shortest cycle by $n$; clearly, $\phi^{n}$ is not the identity but has at least one fixed point. But the set of all fixed points of $\phi^{n}$ is a proper subalgebra of $\mathfrak{A}$ and therefore in class five of $R B L$. Hence our algebra fulfills the hypotheses and is primal.

We will need a couple of rather basic lemmas on subdirect products; a good standard reference with more details on the subject is [1].
2.4.3 Lemma. An algebra $\mathfrak{A}$ is isomorphic to a subdirect product of the algebras $\left\{\mathfrak{A}_{i} \mid i \in I\right\}$ iff for each $i \in I$ there is a homomorphism $h_{i}$ from $\mathfrak{A}$ onto $\mathfrak{A}_{i}$ such that $\bigwedge_{i \in I} \operatorname{ker}\left(h_{i}\right)=0$

Proof. For one implication, consider as homomorphisms the projections $\pi_{i}$ of elements of the subdirect product onto the $i$-th coordinate. The $\pi_{i}$ are obviously homomorphisms and one can easily verify the assertion on the kernels. Conversely, consider the mapping

$$
\phi: \begin{array}{lll}
\mathfrak{A} & \rightarrow & \prod_{i \in I} \mathfrak{A}_{i} \\
a & \mapsto & \left(h_{i}(a)\right)_{i \in I}
\end{array}
$$

Then $\phi$ is clearly a homomorphism onto a subdirect product of $\left\{\mathfrak{A}_{i} \mid i \in I\right\}$ which is one-one and hence an isomorphism as $\bigwedge_{i \in I} \operatorname{ker}\left(h_{i}\right)=0$.
2.4.4 Lemma. Let $\mathfrak{A}$ and $\mathfrak{B}$ be algebras of the same type and let $h: \mathfrak{A} \rightarrow \mathfrak{B}$ be a homomorphism. Then $\operatorname{Con}(h(\mathfrak{A}))$ is isomorphic to the sublattice of $\operatorname{Con}(\mathfrak{A})$ over the set $\{\theta \in \operatorname{Con}(\mathfrak{A}) \mid \theta \geq \operatorname{ker}(h)\}$.

Proof. By the Homomorphism Theorem, $\mathfrak{A} / \operatorname{ker}(h) \cong h(\mathfrak{A})$ so that $\operatorname{Con}(\mathfrak{A} / \operatorname{ker}(h)) \cong$ $\operatorname{Con}(h(\mathfrak{A}))$. But as one can easily see, the sublattice of $\operatorname{Con}(\mathfrak{A})$ over the set $\{\theta \in \operatorname{Con}(\mathfrak{A}) \mid \theta \geq$ $\operatorname{ker}(h)\}$ is isomorphic to $\operatorname{Con}(\mathfrak{A} / \operatorname{ker}(h))$ and the lemma follows.
2.4.5 Lemma. Let $\theta_{1}, \theta_{2}$ be permutable congruences on an algebra $\mathfrak{A}$ satisfying $\theta_{1} \cap \theta_{2}=0$ and $\theta_{1} \cup \theta_{2}=1$. Then $\mathfrak{A} \cong \mathfrak{A} / \theta_{1} \times \mathfrak{A} / \theta_{2}$.

Proof. Consider the mapping

$$
\phi: \begin{array}{ccc}
\mathfrak{A} & \rightarrow & \mathfrak{A} / \theta_{1} \times \mathfrak{A} / \theta_{2} \\
a & \mapsto & \left([a]_{\theta_{1}},[a]_{\theta_{2}}\right)
\end{array}
$$

Apparently, $\phi$ is a homomorphism and since $\theta_{1} \cap \theta_{2}=0$ it is one-one. Now let any $\left([a]_{\theta_{1}},[b]_{\theta_{2}}\right) \in$ $\mathfrak{A} / \theta_{1} \times \mathfrak{A} / \theta_{2}$ be given. As a consequence of the permutability of the two congruences, $\theta_{1} \cdot \theta_{2}=$ $\theta_{1} \cup \theta_{2}=1$ so that $(a, b) \in \theta_{1} \cdot \theta_{2}$. Hence, there exists $z$ such that $(a, z) \in \theta_{1}$ and $(z, b) \in \theta_{2}$. Thus,

$$
\left([a]_{\theta_{1}},[b]_{\theta_{2}}\right)=\left([z]_{\theta_{1}},[z]_{\theta_{2}}\right)=\phi(z)
$$

and $\phi$ is onto.
2.4.6 Lemma. An algebra $\mathfrak{A}$ is isomorphic to the direct product of algebras $\left\{\mathfrak{A}_{1}, \ldots, \mathfrak{A}_{n}\right\}$ iff

1. for each $1 \leq i \leq n$ there is a homomorphism $h_{i}$ from $\mathfrak{A}$ onto $\mathfrak{A}_{i}$ such that
2. the set of all intersections of the kernels of the $h_{i}$ consists of pairwise permutable congruence relations and
3. $\operatorname{ker}\left(h_{1}\right) \cap \ldots \cap \operatorname{ker}\left(h_{n}\right)=0$ and for $2 \leq i \leq n$, $\left(\operatorname{ker}\left(h_{1}\right) \cap \ldots \cap \operatorname{ker}\left(h_{i-1}\right)\right) \cup \operatorname{ker}\left(h_{i}\right)=1$

Proof. If $\mathfrak{A} \cong \mathfrak{A}_{1} \times \ldots \times \mathfrak{A}_{n}$, consider as homomorphisms again the projections; the asserted properties of their kernels are easy to verify. Conversely, by the previous lemma we have that $\mathfrak{A} \cong \mathfrak{B}_{n} \times \mathfrak{A}_{n}$, where $\mathfrak{B}_{n}=\mathfrak{A} /\left(\theta_{1} \cap \ldots \cap \theta_{n-1}\right)$. A straightforward induction finally shows that $\mathfrak{B}_{n} \cong \mathfrak{A}_{1} \times \ldots \times \mathfrak{A}_{n-1}$ and completes the proof.
2.4.7 Theorem. Let $\mathfrak{A}$ be an algebra isomorphic to a subdirect product of finitely many simple algebras $\mathfrak{A}_{1}, \ldots, \mathfrak{A}_{n}$. If the congruences of $\mathfrak{A}$ permute, then $\mathfrak{A}$ is isomorphic to the direct product of a subset of the $\mathfrak{A}_{1}, \ldots \mathfrak{A}_{n}$.

Proof. Let $h_{1}, \ldots, h_{n}$ be the homomorphisms given by Lemma 2.4.3 and denote their kernels by $\theta_{1}, \ldots, \theta_{n}$. Lemma 2.4.4 together with the simplicity of the $\mathfrak{A}_{i}$ implies that the $\theta_{i}$ are maximal in $\operatorname{Con}(\mathfrak{A})$ (if we assume all the $\mathfrak{A}_{i}$ are non-trivial; if not, we simply leave the trivial ones away; if all are trivial then the theorem is as well). Since $\theta_{1} \cap \ldots \cap \theta_{n}=0$, we can extract a minimal subset of the $\theta_{i}$ having the same property. Assume without loss of generality the first $k$ congruences form such a subset, that is, $\theta_{1} \cap \ldots \cap \theta_{k}=0$. Trivially,

$$
\left(\theta_{1} \cap \ldots \cap \theta_{i-1}\right) \cup \theta_{i} \geq \theta_{i}, \quad 2 \leq i \leq k
$$

and thus by the maximality of $\theta_{i},\left(\theta_{1} \cap \ldots \cap \theta_{i-1}\right) \cup \theta_{i}$ must be equal to $\theta_{i}$ or to 1 . But if it was equal to $\theta_{i}$, we could conclude that $\left(\theta_{1} \cap \ldots \cap \theta_{i-1}\right) \leq \theta_{i}$ and then

$$
\theta_{1} \cap \ldots \cap \theta_{i-1} \cap \theta_{i+1} \cap \ldots \cap \theta_{k}=0
$$

contradicting the minimality of the set $\left\{\theta_{1}, \ldots, \theta_{k}\right\}$. Consequently, $\left(\theta_{1} \cap \ldots \cap \theta_{i-1}\right) \cup \theta_{i}=1$ for $2 \leq i \leq k$ and so Lemma 2.4.6 completes the proof.

The following lemma from lattice theory will help us using the congruence distributivity of $\mathcal{V}(\mathfrak{A})$ in our proof. A meet in a lattice is irredundant iff it cannot be written as a meet of a subset of its elements. An element in a lattice is meet irreducible iff it is not the meet of two elements not equal to itself.
2.4.8 Lemma. In a distributive lattice $\mathfrak{L}$, the representation of an element as an irredundant meet of meet-irreducible elements is unique (and dually).

Proof. Let $a$ be an element of $\mathfrak{L}$ such that

$$
a=x_{1} \cap \ldots \cap x_{r}=y_{1} \cap \ldots \cap y_{s}
$$

Then for a fixed $x_{i}$ we can observe the following: Clearly, $x_{i} \geq y_{1} \cap \ldots \cap y_{s}$. Thus $x_{i}=$ $x_{i} \cup\left(y_{1} \cap \ldots \cap y_{s}\right)=\left(x_{i} \cup y_{1}\right) \cap \ldots \cap\left(x_{i} \cup y_{s}\right)$ by the distributivity of $\mathfrak{L}$. But if $x_{i}$ is meet irreducible, the above representation yields $x_{i} \cup y_{j}=x_{i}$ and hence $y_{j} \leq x_{i}$ for some $j$. Similarly for $y_{j}$ we have $x_{k} \leq y_{j}$ for some $k$ so that $x_{k} \leq x_{i}$ and therefore $x_{k}=x_{i}=y_{j}$ because the representation was assumed to be non-redundant. So the $x_{i}$ and $y_{j}$ are equal in pairs, $r=s$ and the representation is indeed unique.

We shall now obtain some results concerning the structure of free algebras. Let $S$ be an algebra type. For a set of $S$-identities $\Psi$ we let $\mathfrak{F}_{k}(\Psi)$ denote the free algebra with $k$ generators determined by $\Psi$. If $\mathfrak{A}$ is an algebra then $\Sigma(\mathfrak{A})$ will denote the equations satisfied by $\mathfrak{A}$. Finally, $\mathfrak{F}_{k}(\mathfrak{A})$ is short for $\mathfrak{F}_{k}(\Sigma(\mathfrak{A}))$. Recall the following important fact:
2.4.9 Lemma. $\mathfrak{F}_{k}(\mathfrak{A}) \in \mathcal{V}(\mathfrak{A})$.

Let $\mathfrak{A}$ be a non-trivial finite algebra of order $n$ and let $G=\left\{\xi_{1}, \ldots, \xi_{k}\right\}$ be a set of $k$ indeterminates. Then there exist $n^{k}$ functions $e_{1}, \ldots, e_{n^{k}}$ from $G$ to $\mathfrak{A}$. Each of the $e_{i}$ induces a subuniverse $S_{i}$ of $\mathfrak{A}$; but as we assume that $\mathfrak{A}$ has no proper subalgebras, all of the $S_{i}$ are equal to $A$.

Construct $\mathfrak{F}_{k}(\mathfrak{A})$ over $G$. All the $e_{i}$ induce in a canonical way a homomorphism $h_{i}$ from $\mathfrak{F}_{k}(\mathfrak{A})$ onto $\mathfrak{A}$ : For a class $\Phi$ of equivalent expressions in $\mathfrak{F}_{k}(\mathfrak{A})$,

$$
h_{i}(\Phi)=\phi\left(e_{i}\left(\xi_{1}\right), \ldots, e_{i}\left(\xi_{n}\right)\right),
$$

where $\phi$ is an arbitrary $S$-expression in $\Phi$. The function is well-defined since for $\phi_{1}, \phi_{2}$ in $\Phi$, $\phi_{1}\left(e_{i}\left(\xi_{1}\right), \ldots, e_{i}\left(\xi_{n}\right)\right)=\phi_{2}\left(e_{i}\left(\xi_{1}\right), \ldots, e_{i}\left(\xi_{k}\right)\right)$. It is easily seen that $h_{i}$ is indeed a homomorphism onto $\mathfrak{A}$. Hence, $\mathfrak{F}_{k}(\mathfrak{A}) / \operatorname{ker}\left(h_{i}\right) \cong \mathfrak{A}$.

If $\Phi_{1} \equiv \Phi_{2}\left(\bigwedge_{1 \leq i \leq n^{k}} \operatorname{ker}\left(h_{i}\right)\right)$ then for all $1 \leq i \leq n^{k}$ we have that $h_{i}\left(\Phi_{1}\right)=h_{i}\left(\Phi_{2}\right)$ and hence for all $\phi_{1} \in \Phi_{1}$ and all $\phi_{2} \in \Phi_{2}, \phi_{1}\left(e_{i}\left(\xi_{1}\right), \ldots, e_{i}\left(\xi_{n}\right)\right)=\phi_{2}\left(e_{i}\left(\xi_{1}\right), \ldots, e_{i}\left(\xi_{k}\right)\right)$. Since this holds for all possible mappings $e_{i}$ into $\mathfrak{A}, \phi_{1}=\phi_{2}$ must be an identity of $\Sigma(\mathfrak{A})$ and thus $\Phi_{1}=\Phi_{2}$. Therefore,

$$
\begin{equation*}
\bigwedge_{1 \leq i \leq n^{k}} \operatorname{ker}\left(h_{i}\right)=0 \tag{2.4.1}
\end{equation*}
$$

and using Lemma 2.4.3 we conclude:
2.4.10 Lemma. Let $\mathfrak{A}$ be a non-trivial finite algebra of order $n$ having no proper subalgebras. Then $\mathfrak{F}_{k}(\mathfrak{A})$ is isomorphic to a subdirect product of $n^{k}$ copies of $\mathfrak{A}$ via the mapping

$$
\begin{array}{clc}
\mathfrak{F}_{k}(\mathfrak{A}) & \rightarrow & \mathfrak{A}^{n^{k}} \\
\Phi & \mapsto & {\left[\phi\left(e_{1}\left(\xi_{1}\right), \ldots, e_{1}\left(\xi_{k}\right)\right), \ldots, \phi\left(e_{n^{k}}\left(\xi_{1}\right), \ldots, e_{n^{k}}\left(\xi_{k}\right)\right)\right]}
\end{array}
$$

where $\phi$ is an arbitrary term in $\Phi$.

With the additional assumption that $\mathfrak{A}$ is simple and generates a congruence permutable variety, Theorem 2.4.7 and Lemma 2.4.9 imply
2.4.11 Lemma. Let $\mathfrak{A}$ be a simple non-trivial finite algebra of order $n$ having no proper subalgebras, and which generates a congruence permutable equational class. Then there exists $1 \leq r \leq n^{k}$ such that $\mathfrak{F}_{k}(\mathfrak{A})$ is isomorphic to $\mathfrak{A}^{r}$.

If all of the assumptions on $\mathfrak{A}$ in Theorem 2.4.1 hold, then all of the factors occur in the representation of the free algebra $\mathfrak{F}_{k}(\mathfrak{A})$.
2.4.12 Lemma. If $\mathfrak{A}$ is an algebra satisfying the assumptions of Theorem 2.4.1, then $\mathfrak{F}_{k}(\mathfrak{A}) \cong$ $\mathfrak{A}^{n^{k}}$.

Proof. First note that the kernels $\operatorname{ker}\left(h_{i}\right)$ must be distinct. For assume that $\operatorname{ker}\left(h_{j}\right)=\operatorname{ker}\left(h_{i}\right)$ for some $j \neq i$. Then

$$
\mathfrak{A} \cong \mathfrak{F}_{k}(\mathfrak{A}) / \operatorname{ker}\left(h_{j}\right)=\mathfrak{F}_{k}(\mathfrak{A}) / \operatorname{ker}\left(h_{i}\right) \cong \mathfrak{A}
$$

and so, since $h_{i}$ and $h_{j}$ are different homomorphisms, we have found a non-trivial automorphism on $\mathfrak{A}$ contrary to our assumption. Now since the kernels are maximal by Lemma 2.4.4, they are meet irreducible. Therefore, equation (2.4.1) provides a representation of 0 as a meet of meet irreducible elements. Since $\mathfrak{A}$ generates a congruence distributive variety and since $\mathfrak{F}_{k}(\mathfrak{A}) \in \mathcal{V}(\mathfrak{A})$, the congruence lattice of $\mathfrak{F}_{k}(\mathfrak{A})$ is distributive. Now assume the representation of 0 as the meet of the kernels of our homomorphisms can be shortened; say

$$
\bigwedge_{\substack{1 \leq l \leq n^{k} \\ l \neq i}} \operatorname{ker}\left(h_{l}\right)=\bigwedge_{1 \leq l \leq n^{k}} \operatorname{ker}\left(h_{l}\right)=0
$$

for some $1 \leq i \leq n^{k}$. Then, just like in the proof of Theorem 2.4.8, we get that there exists $j \neq i$ such that $\operatorname{ker}\left(h_{i}\right) \geq \operatorname{ker}\left(h_{j}\right)$, contradicting either the maximality or the distinctiveness of the kernels.

Now we can prove $\mathfrak{A}$ primal.
Proof of Theorem 2.4.1. Let $f$ be any $k$-ary function on $A$. Then by the preceding lemma, there exists a class $\Phi$ in $\mathfrak{F}_{k}(\mathfrak{A})$ such that for every $\phi \in \Phi$ and for all $1 \leq i \leq n^{k}$ the identity

$$
\phi\left(e_{i}\left(\xi_{1}\right), \ldots, e_{i}\left(\xi_{k}\right)\right)=f\left(e_{i}\left(\xi_{1}\right), \ldots, e_{i}\left(\xi_{k}\right)\right)
$$

holds. But since $\left(e_{i}\left(\xi_{1}\right), \ldots, e_{i}\left(\xi_{k}\right)\right)$ runs through all $k$-tuples of elements of $\mathfrak{A}$, the term operation $\phi$ is identical with $f$ and so $\mathfrak{A}$ is primal.

## Chapter 3

## The clones from $R B L$ are maximal

In the previous chapter we demonstrated that all maximal clones over a finite set are of the form $\operatorname{Pol}(\rho)$, where $\rho$ is a relation in $R B L$, and provided a characterization of primal algebras. This chapter is devoted to the proof of the converse statement, namely that all clones of that kind are indeed maximal. It will therefore result in the aim of this work, the characterization of maximal clones or preprimal algebras respectively. We will consider the six types of relations of our main theorem one after the other. For the first three of them, the same technique (see M. Goldstern and S. Shelah [7]) will be used for the proof. Each of the other three requires a special treatment; in those cases, we will essentially follow the original proof of I. G. Rosenberg and include a result of J. Słupecki on functional completeness.
3.0.13 Definition. For a set of functions $F \subseteq \mathcal{F}$, we define the closure $\langle F\rangle$ of $F$ to be the smallest clone containing $F$.

It is thus our goal to prove $<\operatorname{Pol}(\rho) \cup\{g\}>=\mathcal{F}$ for all relations $\rho$ in $R B L$ and all $g \notin \operatorname{Pol}(\rho)$.
3.0.14 Notation. Throughout this chapter, the letter $\kappa$ will be reserved to denote the cardinality of our finite base set $A$. Moreover, $\left\{\alpha_{1}, \ldots, \alpha_{\kappa}\right\}=A$ will be a fixed enumeration of $A$.

### 3.1 Partial orders with least and greatest element

Let $\rho \subseteq A^{2}$ be a partial order with least and greatest element. For $a, b \in A^{n}$ we write $a \leq b$ iff $\left(a_{i}, b_{i}\right) \in \rho$ for all $1 \leq i \leq n$ and $a<b$ iff $a \leq b$ and $a \neq b$. Let $g \notin \operatorname{Pol}(\rho)$ be an $n$-ary non-monotone function, that is there exist $a, b \in A^{n}$ such that $a \leq b$ but $g(a) \not \leq g(b)$. Since $\rho$ has a greatest element (and $|A|>1$ ), it is non-trivial and $g$ exists.
3.1.1 Theorem. If $\rho \subseteq A^{2}$ is a partial order with least and greatest element, then $\operatorname{Pol}(\rho)$ is a maximal clone.
3.1.2 Lemma. For any $k$ and all $c, d \in A^{k}, c<d$, there exists $f_{c d} \in\langle\operatorname{Pol}(\rho) \cup\{g\}>$ such that $f_{c d}(c) \not \leq f_{c d}(d)$.

Proof. Our first step is to see that for given $c, d \in A, c<d$, we can construct an unary $f_{c d} \in<\operatorname{Pol}(\rho) \cup\{g\}>$ satisfying $f_{c d}(c) \not \leq f_{c d}(d)$ : There are unary monotone functions $f_{c d}^{i}(x)$ mapping $c$ to $a_{i}$ and $d$ to $b_{i}, 1 \leq i \leq n$. This is because we can map all elements $s$ with $s \leq c$ to $a_{i}$, and all other elements to $b_{i}$. Now if $f_{c d}^{i}(y) \not \leq f_{c d}^{i}(z)$, then $f_{c d}^{i}(z)=a_{i}$ which implies $z \leq c$. But as $f_{c d}^{i}(y)$ must equal $b_{i}, y \not \leq c$ and so $y \not \leq z$. Hence, the functions $f_{c d}^{i}(x)$ are indeed monotone. Now set $f_{c d}(x)=g\left(f_{c d}^{1}(x), \ldots, f_{c d}^{n}(x)\right)$. Then $f_{c d}(c)=g(a) \not \leq g(b)=f_{c d}(d)$ which is exactly what we wanted.
Next note that we can do the same thing for arbitrary tuples $c, d \in A^{k}$ with $c<d$ : Choose $1 \leq i \leq n$ such that $c_{i} \neq d_{i}$. Since $c_{i}<d_{i}$, we can construct $f_{c_{i} d_{i}}$ as shown before and then define $f_{c d}=f_{c_{i} d_{i}} \circ \pi_{i}^{k}$.

Proof of Theorem 3.1.1. Let $h$ be an arbitrary $k$-ary function. Using the functions we just constructed in the preceding lemma, we will show $h \in\langle\operatorname{Pol}(\rho) \cup\{g\}\rangle$. Consider the set $S=\left\{f_{c d} \mid c, d \in A^{k}, c<d\right\}$, denote it for reasons of simpler notation by $\left\{f_{i} \mid i \in I\right\}$, and define a mapping

$$
\text { ext: } \begin{array}{cccc}
A^{k} & \rightarrow & A^{k+|I|} \\
x & \mapsto & \left(x,\left(f_{i}(x)\right)_{i \in I}\right) .
\end{array}
$$

Then for all distinct $x, y \in A^{k}$ we have that $\operatorname{ext}(x) \not \leq \operatorname{ext}(y)$. This is trivial if $x \not \leq y$, and if otherwise, then the function $f_{x y}$ satisfying $f_{x y}(x) \not \leq f_{x y}(y)$ is an element of $S$ so that by the definition of $\operatorname{ext}, \operatorname{ext}(x) \not \leq \operatorname{ext}(y)$. Now define an operation $H$ on the range $\left\{\operatorname{ext}(x) \mid x \in A^{k}\right\}$ of ext by $H(\operatorname{ext}(x))=h(x)$. $H$ respects $\rho$ as on its domain no elements are comparable. We can find a monotone continuation $\tilde{H}$ of $H$ by setting for all $x$ not in the range of ext

$$
\tilde{H}(x)= \begin{cases}o & , \exists y \in A^{k}(\operatorname{ext}(y) \leq x) \\ z & , \text { otherwise }\end{cases}
$$

where $o$ is the greatest element and $z$ the least element of $\rho$. But $\tilde{H} \in \operatorname{Pol}(\rho)$, and so, as obviously $h(x)=\tilde{H}\left(x,\left(f_{i}(x)\right)_{i \in I}\right)$, we get that $h \in<\operatorname{Pol}(\rho) \cup\{g\}>$.

### 3.2 Non-trivial equivalence relations

Let $\rho \subseteq A^{2}$ be a non-trivial equivalence relation on $A$. For $a, b \in A^{n}$ we write $a \sim b$ iff $\left(a_{i}, b_{i}\right) \in \rho$ for all $1 \leq i \leq n$. Obviously $\sim$ is an equivalence relation on $A^{n}$. Let $g \notin \operatorname{Pol}(\rho)$ be an $n$-ary function not preserving $\rho$, that is, there are $a, b \in A^{n}$ such that $a \sim b$ but $g(a) \nsim g(b)$. As $\rho$ is non-trivial, $g$ exists.
3.2.1 Theorem. If $\rho \subseteq A^{2}$ is a non-trivial equivalence relation, then $\operatorname{Pol}(\rho)$ is a maximal clone.

Just like with partial orders, the following lemma is a fact.
3.2.2 Lemma. For any $k$ and all distinct $c, d \in A^{k}, c \sim d$, there exists $f_{c d} \in<\operatorname{Pol}(\rho) \cup\{g\}>$ such that $f_{c d}(c) \nsim f_{c d}(d)$.

Proof. As in Lemma 3.1.2, for arbitrary distinct $c, d \in A, c \sim d$, we construct an unary $f_{c d} \in<\operatorname{Pol}(\rho) \cup\{g\}>$ such that $f_{c d}(c) \nsim f_{c d}(d)$. Define functions $f_{c d}^{i}(x), i=1, \ldots, n$ by mapping $c$ to $a_{i}$ and all other elements to $b_{i}$. Obviously, as $a_{i} \sim b_{i}, f_{c d}^{i} \in \operatorname{Pol}(\rho)$ for all $1 \leq i \leq n$, and setting $f_{c d}(x)=g\left(f_{c d}^{1}(x), \ldots, f_{c d}^{n}(x)\right)$ yields the desired function.
For arbitrary distinct tuples $c, d \in A^{k}$ with $c \sim d$, we define $f_{c d}=f_{c_{i} d_{i}} \circ \pi_{i}^{k}$, where $i$ is arbitrary with $c_{i} \neq d_{i}$, and the lemma follows.

Proof of Theorem 3.2.1. Let $h$ be an arbitrary $k$-ary function. Following the proof of Theorem 3.1.1, we define the functions ext and $H$. Again, since no elements in the image of ext are equivalent with respect to $\sim$, we can extend $H$ to $\tilde{H} \in \operatorname{Pol}(\rho)$ by mapping all members of an equivalence class $e$ to a fixed element $x_{e}$ of $A$. The element $x_{e}$ is determined if $\operatorname{ext}(x) \in e$ for some $x \in A^{k}$; otherwise, it can be chosen arbitrarily. Hence, $h(x)=\tilde{H}\left(x,\left(f_{i}(x)\right)_{i \in I}\right) \in$ $<\operatorname{Pol}(\rho) \cup\{g\}>$.

### 3.3 Prime permutations

Let $\rho \subseteq A^{2}$ be the graph of a prime permutation $\pi$ on $A$. For an element $a$ of $A^{n}$ and $l \geq 1$ we write $a+l$ for the $n$-tuple $\left(\pi^{l}\left(a_{1}\right), \ldots, \pi^{l}\left(a_{n}\right)\right)$. Then on $A,(a, b) \in \rho$ means exactly that $a+1=b$. We call two elements $a, b \in A^{n}$ parallel iff there is an $l \geq 1$ such that $a+l=b$. Clearly, by that notion an equivalence relation is defined on $A^{n}$ for every $n$. Let $g \notin \operatorname{Pol}(\rho)$ be an $n$-ary function not preserving $\rho$, that is, there are $a, b \in A^{n}$ such that $a+1=b$ but $g(a)+1 \neq g(b)$.
3.3.1 Theorem. If $\rho \subseteq A^{2}$ is a prime permutation, then $\operatorname{Pol}(\rho)$ is a maximal clone.

Similarly to the preceding two cases we have:
3.3.2 Lemma. Let $k \geq 1$ and $c \in A^{k}$ with $c+l=d$ for some $1 \leq l \leq p-1$. Then there exists $f_{c d} \in<\operatorname{Pol}(\rho) \cup\{g\}>$ such that $f_{c d}(c)+l \neq f_{c d}(d)$.

Proof. Our first assertion is that there are $\tilde{a}, \tilde{b} \in A^{n}$ with $\tilde{a}+l=\tilde{b}$ but $g(\tilde{a})+l \neq g(\tilde{b})$. For assume $g(\tilde{a}+l)=g(\tilde{a})+l$ for all $\tilde{a} \in A^{n}$; then if we add $l$ to $a$ for $l^{-1}$ times, where $l^{-1}$ is the multiplicative inverse of $l$ modulo $p$, we get that $g(a+1)=g\left(a+l^{-1} l\right)=g(a)+l^{-1} l=$ $g(a)+1$, contradiction. Now if $c, d \in A$, it is clear that there are functions $f_{c d}^{i} \in \operatorname{Pol}(\rho)$ such that $f_{c d}^{i}(c)=\tilde{a}_{i}$ and $f_{c d}^{i}(d)=f_{c d}^{i}(c+l)=f_{c d}^{i}(c)+l=\tilde{b}_{i}$ for all $1 \leq i \leq n$. Thus,
$f_{c d}(x)=g\left(f_{c d}^{1}(x), \ldots, f_{c d}^{n}(x)\right) \in<\operatorname{Pol}(\rho) \cup\{g\}>$ satisfies the assertion of the lemma. In the case of tuples $c, d \in A^{k}$, we do as before and set $f_{c d}=f_{c_{1} d_{1}} \circ \pi_{1}^{k}$.

Proof of Theorem 3.2.1. Let $h$ be an arbitrary $k$-ary function. Again we define the functions ext and $H$. Now obviously no elements in the image of ext are parallel. Since the value of an element under a function in $\operatorname{Pol}(\rho)$ determines only the values of its parallel class, we find an extension $\tilde{H}$ of $H$ such that $\tilde{H} \in \operatorname{Pol}(\rho)$. Therefore, as $h(x)=\tilde{H}\left(x,\left(f_{i}(x)\right)_{i \in I}\right) \in\langle\operatorname{Pol}(\rho) \cup\{g\}>$, it follows that $\operatorname{Pol}(\rho)$ is a maximal clone.

### 3.4 Central relations

We will show that every central relation $\rho \subseteq A^{h}$ yields a maximal clone via Pol. We distinguish the possibilities $h=1$, in which case $\rho$ is just a proper subset of $A$, and $h \geq 2$. In the first case, the method we used so far can be applied once again; however, in all other cases the issue is more complicated. As before we denote by $g \notin \operatorname{Pol}(\rho)$ the $n$-ary function not preserving $\rho ; g$ exists as the center of a central relation is non-trivial by definition. Thus, there exist $a_{1}, \ldots, a_{n} \in \rho$ such that $\left(g\left(a_{11}, \ldots, a_{n 1}\right), \ldots, g\left(a_{1 h}, \ldots, a_{n h}\right)\right) \notin \rho$. The following theorem does the case $h=1$.
3.4.1 Theorem. If $\rho \subseteq A$ is a proper subset of $A$, then $\operatorname{Pol}(\rho)$ is a maximal clone.
3.4.2 Lemma. For every $c \in \rho^{k}$ there is a $f_{c} \in\left\langle\operatorname{Pol}(\rho) \cup\{g\}>\right.$ with $f_{c}(c) \notin \rho$.

Proof. There are $a_{1}, \ldots, a_{n} \in \rho$ such that $g\left(a_{1}, \ldots, a_{n}\right) \notin \rho$. If $c \in \rho$, then there are obviously mappings $f_{c}^{i} \in \operatorname{Pol}(\rho)$ with $f_{c}^{i}(c)=a_{i}$. Setting $f_{c}=g\left(f_{c}^{1}, \ldots, f_{c}^{n}\right)$ proves the lemma for this case. If $c$ is a $k$-tuple, define $f_{c}=f_{c_{1}} \circ \pi_{1}^{k}$ as usually.

Proof of Theorem 3.4.1. Take any $k$-ary function $h$ and define for every $x \in A^{k}$ the tuple $\operatorname{ext}(x)$ by $\operatorname{ext}(x)=\left(x,\left(f_{c}(x)\right)_{c \in \rho^{k}}\right)$. On the range of ext, set $H(\operatorname{ext}(x))=h(x)$. Clearly, as tuples of the form $\operatorname{ext}(x)$ can never have all their components in $\rho$, we can extend $H$ to $\tilde{H} \in \operatorname{Pol}(\rho)$ like in the previous sections and as $h(x)=\tilde{H}\left(x,\left(f_{c}(x)\right)_{c \in \rho^{k}}\right)$ the theorem has been proven.

## A completeness criterion

For the central relations as well as the $h$-regularly generated relations we will need a completeness criterion due to J. Slupecki saying that for $|A| \geq 3$, if $F \subseteq \mathcal{F}$ contains all unary functions and a function which takes all values of $A$ and which depends on at least two variables, then $<F>=\mathcal{F}$. This criterion will be proven now; we will essentially follow a proof by J. W. Butler in [3]. The restriction $|A| \geq 3$ does not matter to us: For central relations we use the criterion only for the case $2 \leq h<\kappa$; in the case of $h$-regularly generated relations, $3 \leq h \leq \kappa$ by definition.
3.4.3 Definition. An $n$-ary function $f\left(x_{1}, \ldots, x_{n}\right)$ depends on the $j$-th variable, $1 \leq j \leq n$, iff there exist $a \in A^{n}$ and $u \in A$ such that $f\left(a_{1}, \ldots, a_{n}\right) \neq f\left(a_{1}, \ldots, a_{j-1}, u, a_{j+1}, \ldots, a_{n}\right)$. We call $f$ irreducible iff it depends on at least two variables and reducible iff it does not.

We denote the range of a function $f$ by $\Re(f)$.
3.4.4 Lemma. Let $\kappa \geq 3$ and let $f$ be an irreducible function of $n$ arguments, $n \geq 3$, which is onto. Then there is an irreducible function $g$ of two variables in $<\mathcal{F}_{1} \cup\{f\}>$ which is onto.

Proof. There are $1 \leq q \leq n, a \in A^{n}$ and $u \in A$ such that $f(\tilde{a}) \neq f(a)$ if we set $\tilde{a}=$ $\left(a_{1}, \ldots, a_{q-1}, u, a_{q+1}, \ldots, a_{n}\right)$. Say $f(a)=\alpha_{1}$ and $f(\tilde{a})=\alpha_{2}$ and choose $y_{i} \in A^{n}, 3 \leq i \leq \kappa$ such that $f\left(y_{i}\right)=\alpha_{i}$. Note next that there exist $w, z \in A^{n}$ with $w_{q}=z_{q}$ but $f(w) \neq f(z)$, for otherwise $f$ would depend only on its $q$-th argument and would therefore be reducible. We distinguish two cases: First, such $w$ and $z$ exist with the additional property that $f(w) \neq f(a)$ and $f(w) \neq f(\tilde{a})$, and second, no such $w$ and $z$ fulfill this additional assumption.
In the first case, say without loss of generality $f(w)=f\left(y_{3}\right)$. We define $n-1$ unary functions $h_{i}, 1 \leq i \leq n, i \neq q$ by

$$
h_{i}(x)= \begin{cases}a_{i} & , x=\alpha_{1} \\ z_{i} & , x=\alpha_{2} \\ w_{i} & , x=\alpha_{3} \\ y_{j i} & , x=\alpha_{j} \wedge j \notin\{1,2,3\}\end{cases}
$$

and $h_{q}$ by

$$
h_{q}(x)= \begin{cases}a_{q} & , x=\alpha_{1} \\ u & , x=\alpha_{2} \\ z_{q} & , x=\alpha_{3} \\ y_{j q} & , x=\alpha_{j} \wedge j \notin\{1,2,3\}\end{cases}
$$

and set $g(x, y)=f\left(h_{1}(x), \ldots, h_{q-1}(x), h_{q}(y), h_{q+1}(x), \ldots, h_{n}(x)\right)$. Then $g$ is onto since $g\left(\alpha_{1}, \alpha_{1}\right)=f(a)=\alpha_{1}, g\left(\alpha_{1}, \alpha_{2}\right)=f(\tilde{a})=\alpha_{2}, g\left(\alpha_{3}, \alpha_{3}\right)=f(w)=\alpha_{3}$, and $g\left(\alpha_{i}, \alpha_{i}\right)=f\left(y_{i}\right)=$ $\alpha_{i}$ for $i>3$. Moreover, $g$ is not reducible as $g\left(\alpha_{1}, \alpha_{1}\right) \neq g\left(\alpha_{1}, \alpha_{2}\right)$ and $g\left(\alpha_{2}, \alpha_{3}\right) \neq g\left(\alpha_{3}, \alpha_{3}\right)$. In the second case, we choose for $1 \leq i \leq n, i \neq q$ functions $h_{i}$ satisfying

$$
h_{i}(x)= \begin{cases}a_{i} & , x=\alpha_{1} \\ w_{i} & , x=\alpha_{2} \\ z_{i} & , x=\alpha_{3}\end{cases}
$$

Define $g(x, y)=f\left(h_{1}(x), \ldots, h_{q-1}(x), y, h_{q+1}(x), \ldots, h_{n}(x)\right)$. Now the condition of this case implies that if $s \in A^{n}$ and $s_{q}=y_{i q}$ for some $3 \leq i \leq \kappa$, then $f(s)=f\left(y_{i}\right)$ since $f\left(y_{i}\right) \neq f(a)$ and $f\left(y_{i}\right) \neq f(\tilde{a})$. Thus, $g\left(\alpha_{1}, a_{q}\right)=g(a)=\alpha_{1}, g\left(\alpha_{1}, u\right)=f(\tilde{a})=\alpha_{2}$, and $g\left(\alpha_{m}, y_{i q}\right)=f\left(y_{i}\right)=\alpha_{i}$ for any $1 \leq m \leq n$ and $3 \leq i \leq \kappa$. Hence, $g$ is onto. Moreover, $g\left(\alpha_{1}, a_{q}\right)=g(a) \neq g(\tilde{a})=$ $g\left(\alpha_{1}, u\right)$ and $g\left(\alpha_{2}, w_{q}\right)=g(w) \neq g(z)=g\left(\alpha_{3}, w_{q}\right)$ and so $g$ is irreducible.
3.4.5 Lemma. If $f \in \mathcal{F}_{2}$ is an irreducible function of two variables which takes at least three distinct values, then there exist $a, b, c, d \in A$ such that $f$ takes three distinct values on $\{(a, c),(a, d),(b, c),(b, d)\}$.

Proof. Assume first that there is an $a \in A$ such that $f$ takes at least three values on $\{(a, x) \mid x \in$ $A\}$. Since $f$ is irreducible, there must be $b, c \in A$ such that $f(a, c) \neq f(b, c)$. As $f$ takes at least three values with $a$ as the first argument, there is $d \in A$ with $f(a, d) \neq f(a, c)$ and $f(a, d) \neq f(b, c)$.
Consider now the case where there is no such $a$. The irreducibility of $f$ implies there is $a \in A$ such that $f$ takes two values with $a$ as the first argument. It follows from the assumption for this case that there is $w$ in the range of $f$ such that $f(a, x) \neq w$ for all $x \in A$; say $w=f(b, c)$, $b \neq a$. Hence, $f(a, c) \neq f(b, c)$. Now take any $d \in A$ with $f(a, c) \neq f(a, d)$ to finish the proof.
3.4.6 Lemma. If $f \in \mathcal{F}_{2}$ is an irreducible function of two variables with $|\Re(f)|=p, p \geq 3$, then there exist two unary functions $h_{1}, h_{2} \in \mathcal{F}_{1}$ which both take at most $p-1$ elements such that for every $x \in \Re(f)$ we have $f\left(h_{1}(x), h_{2}(x)\right)=x$.

Proof. Let $a, b, c, d$ be provided by Lemma 3.4.5. Assume without loss of generality that $f(a, c)=u, f(a, d)=v, f(b, c)=w$ are all different. We define $h_{1}, h_{2} \in \mathcal{F}_{1}$ as follows: For $u, v, w$ we set $h_{1}(u)=a, h_{1}(v)=a, h_{1}(w)=b$ and $h_{2}(u)=c, h_{2}(v)=d, h_{2}(w)=c$; for $x \notin \Re(f)$, we define $h_{1}(x)=a$ and $h_{2}(x)=c$; and for $x \in \Re(f) \backslash\{u, v, w\}$ we choose any values for $h_{1}(x)$ and $h_{2}(x)$ such that the requirement $f\left(h_{1}(x), h_{2}(x)\right)=x$ is satisfied. Clearly, $h_{1}, h_{2}$ have all desired properties.

To proof the completeness criterion, we want to construct the function returning the maximum of to elements with respect to some total ordering of the elements of $A$. Therefore we will for the rest of this section replace $A$ by the set of natural numbers $\kappa=\{0, \ldots, \kappa-1\}$ and use the standard notions of $\leq, \vee,+$ and - on that set.
3.4.7 Lemma. Let $f \in \mathcal{F}_{2}, p<\kappa$, and assume there exist $i, j, l \in A$ such that for all $y<p$ we have $f(i, y)=y$ and $f(j, y)=l$. Then there is a function of two variables $g \in<\mathcal{F}_{1} \cup\{f\}>$ such that $g(x, y)=x \vee y$ for $x, y<p$.

Proof. We may assume without loss of generality that $i, j, l<p$. This is legitimate as we can shift them with unary functions. The proof will be by induction on $p$. First, let $p=2$; then there are four possibilities: either $i=0, j=1, l=1$ or $i=0, j=1, l=0$ or $i=1, j=0, l=0$ or $i=1, j=0, l=1$. In multiplication tables of the restriction of $f$ to $\{0,1\}^{2}$, these scenarios look like this:

(i) |  | 0 | 1 |
| :---: | :---: | :---: |
| 0 | 0 | 1 |
| 1 | 1 | 1 |

(ii) $\left.\begin{array}{c|cc} & & 0 \\ \hline\end{array} \begin{array}{l}1 \\ \hline 0 \\ 1\end{array}\right)$
(iii) $\left.\begin{array}{c|cc} & & 0 \\ \hline\end{array} \begin{array}{l}1 \\ \hline 0 \\ 1\end{array}\right) 0 \begin{aligned} & 0 \\ & \\ & \end{aligned}$
(iv) $\left.\begin{array}{c|cc} & & 0 \\ \hline\end{array} \begin{array}{l}1 \\ \hline 0 \\ 1\end{array}\right)$

In the first case, $f$ is the maximum function on $\{0,1\}$ and we can take $f$ itself for $g$. In the other cases we use any unary function $h$ exchanging 0 and 1 to define $g(x, y)$ to be $h(f(x, h(y)))$ or $h(f(h(x), h(y)))$ or $f(h(x), y)$, respectively.
Now assume the lemma is true for $p-1$, and let $g^{\prime}$ be a function in $<\mathcal{F}_{1} \cup\{f\}>$ satisfying $g^{\prime}(x, y)=x \vee y$ for $x, y<p-1$. Choose functions $h_{1}, h_{2} \in \mathcal{F}_{1}$ such that

$$
h_{1}(x)=\left\{\begin{array}{ll}
i & , x<p-1 \\
j & , x=p-1
\end{array} \quad h_{2}(x)= \begin{cases}p-1 & , x=l \\
l & , x=p-1 \\
x & , \text { otherwise }\end{cases}\right.
$$

and construct $f^{\prime} \in \mathcal{F}_{2}$ as $f^{\prime}(x, y)=h_{2}\left(f\left(h_{1}(x), h_{2}(y)\right)\right)$. It is easy to check that

$$
f^{\prime}(x, y)= \begin{cases}y & , x<p-1 \wedge y<p \\ p-1 & , x=p-1 \wedge y<p\end{cases}
$$

Now we define $g$ by $g(x, y)=f^{\prime}\left(f^{\prime}(x, y), g^{\prime}(x, y)\right)$. One readily verifies that for $x, y<p-1$, $g(x, y)=g^{\prime}(x, y)=x \vee y$; for $x=p-1, y<p, g(x, y)=f^{\prime}\left(p-1, g^{\prime}(p-1, y)\right)=p-1$; and for $x<p-1, y=p-1, g(x, y)=f^{\prime}\left(p-1, g^{\prime}(x, p-1)\right)=p-1$. Hence, $g(x, y)=x \vee y$ for $x, y<p$.
3.4.8 Lemma. If $f \in \mathcal{F}_{2}$ is irreducible and $\Re(f)=\{0, \ldots, p-1\}, 3 \leq p \leq \kappa$, then there is a binary function $g \in<\mathcal{F}_{1} \cup\{f\}>$ such that $g(x, y)=x \vee y$ for $x, y<p$.

Proof. The proof will be by induction on $p$. If $p=3$, by Lemma 3.4.5 there are $a, b, c, d$ such that $f$ takes at least three distinct values on $\{(a, c),(a, d),(b, c),(b, d)\}$. By shifting those elements and their values under $f$ with unary functions, we may assume that $a=c=0$, $b=d=1, f(0,0)=0, f(0,1)=1$, and $f(1,0)=2$. This leaves us essentially with two possible multiplication tables:

> (ii) |  | 0 | 1 |
| :---: | :---: | :---: |
| 0 | 0 | 1 |
| 1 | 2 | 0 |

In the first case, we choose functions $h_{1}, h_{2} \in \mathcal{F}_{1}$ with

$$
\begin{array}{ll}
h_{1}(0)=0 & h_{2}(0)=0 \\
h_{1}(1)=0 & h_{2}(1)=1 \\
h_{1}(2)=1 & h_{2}(2)=1 .
\end{array}
$$

Then we can construct $g$ as

$$
g(x, y)=f\left(h_{2} \circ f\left(h_{1}(x), h_{1}(y)\right), h_{2} \circ f\left(h_{2}(x), h_{2}(y)\right)\right)
$$

To construct $g$ in the other case we choose additional functions $h_{3}, h_{4} \in \mathcal{F}_{1}$ with

$$
\begin{array}{ll}
h_{3}(0)=2 & h_{4}(0)=1 \\
h_{3}(1)=0 & h_{4}(1)=0 \\
h_{3}(2)=1 &
\end{array}
$$

and define $g^{\prime} \in \mathcal{F}_{2}$ by

$$
g^{\prime}(x, y)=h_{3} \circ f\left(y, h_{2} \circ f\left(x, h_{4}(y)\right)\right)
$$

It is boring but possible to verify that $g^{\prime}$ agrees on $\{0,1\}$ with the $f$ of the first case which we already treated.
Assuming our assertion is true for $p-1$, we prove it for $p, 3<p \leq n$. First we construct a function $f^{\prime \prime}$ from $f$ satisfying the hypotheses of the lemma for $p-1$; we need to restrict the range of $f$ to $\{0, \ldots, p-2\}$ without making $f$ reducible. To do this we apply Lemma 3.4.5, taking $a, b, c, d \in A$ such that at least three distinct values $u, v, w \in A$ are represented among $\{f(a, c), f(a, d), f(b, c), f(b, d)\}$. Since $p>3$ there is $z \in \Re(f) \backslash\{u, v, w\}$. Define $h \in \mathcal{F}_{1}$ by

$$
h(x)= \begin{cases}u & , x=z \\ x & , \text { otherwise } .\end{cases}
$$

Then $h(f(x, y)) \in \mathcal{F}_{2}$ is not reducible and has $p-1$ elements in its range. By permuting the elements of $A$ with an unary function we produce a function $f^{\prime \prime} \in \mathcal{F}_{2}$ satisfying the hypotheses of the lemma for $p-1$, and hence by induction hypothesis we get a function $g^{\prime \prime} \in \mathcal{F}_{2}$ such that $g(x, y)=x \vee y$ for $x, y<p-1$.
Next by Lemma 3.4.6 there exist functions $h_{1}, h_{2} \in \mathcal{F}_{1}$ with $\Re\left(h_{1}\right)$, $\Re\left(h_{2}\right)$ consisting of at most $p-1$ elements such that $f\left(h_{1}(x), h_{2}(x)\right)=x$ for $x<p$. There exist permutations $h_{3}, h_{4} \in \mathcal{F}_{1}$ such that $h_{3}(x)<p-1$ for all $x \in \Re\left(h_{1}\right)$ and $h_{4}(x)<p-1$ for all $x \in \Re\left(h_{2}\right)$. Define $h_{5}, h_{6} \in \mathcal{F}_{1}$ and $f^{\prime} \in \mathcal{F}_{2}$ by

$$
\begin{aligned}
h_{5} & =h_{3} \circ h_{1} \\
h_{6} & =h_{4} \circ h_{2} \\
f^{\prime}(x, y) & =f\left(h_{3}^{-1}(x), h_{4}^{-1}(y)\right) .
\end{aligned}
$$

Then obviously $f^{\prime}\left(h_{5}(x), h_{6}(x)\right)=x$ for all $x<p$ and $\Re\left(h_{5}\right), \Re\left(h_{6}\right)$ are subsets of $\{0, \ldots, p-2\}$. We define $g^{\prime} \in \mathcal{F}_{2}$ by $g^{\prime}(x, y)=f^{\prime}\left(g^{\prime \prime}\left(x, h_{5}(y)\right), g^{\prime \prime}\left(x, h_{6}(y)\right)\right)$. Then $g^{\prime}$ satisfies the equation $g^{\prime}(0, y)=f^{\prime}\left(h_{5}(y), h_{6}(y)\right)=y$ for $y<p$; moreover, for $y<p, g^{\prime}(p-2, y)=f^{\prime}(p-2, p-2)$ and is therefore constant. Hence by Lemma 3.4.7 we can generate a function $g \in \mathcal{F}_{2}$ such that $g$ agrees with the maximum function for arguments smaller than $p$.
3.4.9 Theorem. Assume $|A| \geq 3$ and let $f \in \mathcal{F}$ be an irreducible function with $\Re(f)=A$. Then $<\mathcal{F}_{1} \cup\{f\}>=\mathcal{F}$.

Proof. By Lemma 3.4.4 we may assume $f$ is a function of two variables. Lemma 3.4.8 then implies that $<\mathcal{F}_{1} \cup\{f\}>$ contains the maximum function with respect to some total order of
the elements of $A$. But is well-known from the results of E . L. Post and easily verified that the unary functions together with the maximum function already generate all functions of arbitrary arity over $A$; thus, $<\mathcal{F}_{1} \cup\{f\}>=\mathcal{F}$.

## Totally reflexive and totally symmetric relations

The following lemmas hold for totally reflexive and totally symmetric relations. They will help us with both the central relations with $h \geq 2$ and the $h$-regularly generated relations. The first lemma implies that we can assume without loss of generality that the function $g$ not preserving $\rho$ is unary.
3.4.10 Lemma. Let $\rho \neq \iota_{h}^{A}$ be a totally reflexive and totally symmetric $h$-ary relation. If $g \notin \operatorname{Pol}(\rho)$ then there is an unary $f \in<\operatorname{Pol}(\rho) \cup\{g\}>$ that does not preserve $\rho$.

Proof. Let $a_{1}, \ldots, a_{n} \in \rho$ such that $\left(g\left(a_{11}, \ldots, a_{n 1}\right), \ldots, g\left(a_{1 h}, \ldots, a_{n h}\right)\right) \notin \rho$. Choose $\left(c_{1}, \ldots, c_{h}\right) \in \rho,\left(c_{1}, \ldots, c_{h}\right) \notin \iota_{h}^{A}$. Define for $1 \leq i \leq n$ unary functions $f_{i}$ by $f_{i}\left(c_{j}\right)=a_{i j}$, $1 \leq j \leq h$, and $f_{i}(x)=a_{i 1}$ for all other elements $x \in A$. The operations $f_{i}$ preserve $\rho$ as they map just any tuple to a tuple in $\rho$ : If an $h$-tuple consisting of function values of $f_{i}$ has two identical entries, then the tuple is an element of $\rho$ as $\rho$ is totally reflexive; if otherwise, then the definition of $f_{i}$ implies that the tuple contains the values $a_{i 1}, \ldots, a_{i n}$ in some order and is thus in $\rho$ by its total symmetry. Now $f(x)=g\left(f_{1}(x), \ldots, f_{n}(x)\right) \in<\operatorname{Pol}(\rho) \cup\{g\}>$ maps $\left(c_{1}, \ldots, c_{h}\right) \in \rho$ to $\left(g\left(a_{11}, \ldots, a_{n 1}\right), \ldots, g\left(a_{1 h}, \ldots, a_{n h}\right)\right) \notin \rho$.
3.4.11 Lemma. Let $\rho \neq \iota_{h}^{A}$ be a totally reflexive and totally symmetric h-ary relation, $1 \leq$ $h \leq \kappa$. If $g$ is an unary function not preserving $\rho$, then there is a subset $D=\left\{d_{1}, \ldots, d_{h}\right\}$ of A such that $\left(d_{1}, \ldots, d_{h}\right) \notin \rho$ and $<\operatorname{Pol}(\rho) \cup\{g\}>$ contains all unary functions which take only values in $D$.

Proof. There is $\left(a_{1}, \ldots, a_{h}\right) \in \rho$ such that $\left(g\left(a_{1}\right), \ldots, g\left(a_{h}\right)\right) \notin \rho$. If we set $d_{i}=g\left(a_{i}\right), 1 \leq i \leq h$, then $\left(d_{1}, \ldots, d_{h}\right) \notin \rho$. Let $h$ be an unary function that takes only values in $D$. Define a function $l$ by $l(x)=a_{i}$ whenever $h(x)=d_{i}$. Then by the same argument as in the preceding lemma for $f_{i}, l \in \operatorname{Pol}(\rho)$. Hence, $h=g \circ l \in\langle\operatorname{Pol}(\rho) \cup\{g\}>$.
3.4.12 Theorem. Let $\rho \neq \iota_{h}^{A}$ be a totally reflexive and totally symmetric non-trivial $h$-ary relation, where $2 \leq h \leq \kappa$. If for every $D=\left\{d_{1}, \ldots, d_{h}\right\}$ with $\left(d_{1}, \ldots, d_{h}\right) \notin \rho$ an n-ary function $q \in \operatorname{Pol}(\rho)$ exists which takes all values of $A$ on $D^{n}$, then $\operatorname{Pol}(\rho)$ is a maximal clone.

Proof. Take an unary $g \notin \operatorname{Pol}(\rho)$. Let $D$ be provided by Lemma 3.4.11. By our hypothesis, there are $q\left(x_{1}, \ldots, x_{n}\right) \in \operatorname{Pol}(\rho)$ and $a_{i} \in D^{n}, 1 \leq i \leq \kappa$, such that $q\left(a_{i}\right)=\alpha_{i}$. Let $h \in \mathcal{F}_{1}$ be given. Define for $1 \leq j \leq n$ functions $g_{j} \in \mathcal{F}_{1}$ by $g_{j}(x)=a_{i j}$ whenever $h(x)=\alpha_{i}$. As the $g_{j}$ obviously take only values in $D$, we have $g_{j} \in<\operatorname{Pol}(\rho) \cup\{g\}>$ and so the same holds for $q\left(g_{1}, \ldots, g_{n}\right)$. But it is easily verified that $q\left(g_{1}(x), \ldots, g_{n}(x)\right)=h(x)$ for all $x \in A$, and so $h \in<\operatorname{Pol}(\rho) \cup\{g\}>$. We have thus shown that $\langle\operatorname{Pol}(\rho) \cup\{g\}>$ contains all unary functions.

Now assume $q$ depends only on one variable. Then, as $q$ takes all values of $A$, we necessarily have that $D=A$ and so $h=\kappa$. Therefore, $\left(\alpha_{1}, \ldots, \alpha_{\kappa}\right) \notin \rho$. But this implies $\rho=\iota_{h}^{A}$, contradictory to our assumption. Hence, $q$ depends on at least two variables, and we can apply Theorem 3.4.9 to obtain that $\langle\operatorname{Pol}(\rho) \cup\{g\}>=\mathcal{F}$.

## Central relations with $h \geq 2$

We make use of the preceding results for the remaining case $h \geq 2$.
3.4.13 Lemma. Let $\rho$ be an $h$-ary central relation, $2 \leq h \leq \kappa,\left(d_{0}, \ldots, d_{h-1}\right) \notin \rho$, and let $u \in A$ be a central element of $\rho$. Enumerate all functions in $h^{\kappa}$ by $\left\{p_{1}, \ldots, p_{h^{\kappa}}\right\}$. For $i \in \kappa$ set $b_{i}=\left(d_{p_{1}(i)}, \ldots, d_{p_{h^{\kappa}}(i)}\right)$. We define a $h^{\kappa}$-ary function $q$ by $q\left(b_{i}\right)=\alpha_{i+1}, i \in \kappa$, and for all other elements $x \in A^{h^{\kappa}}$ we set $q(x)=u$. Then $q$ preserves $\rho$.

Proof. We first show that for distinct $i_{0}, \ldots, i_{h-1} \in \kappa,\left(b_{i_{0}}, \ldots, b_{i_{h-1}}\right) \notin \rho^{h^{\kappa}}$. Take any function $r \in h^{\kappa}$ with $r\left(i_{j}\right)=j, j \in h$. There is an $1 \leq l \leq h^{\kappa}$ such that $r=p_{l}$. Thus $\left(b_{i_{0} l}, \ldots, b_{i_{h-1} l}\right)=$ $\left(d_{p_{l}\left(i_{0}\right)}, \ldots, d_{p_{l}\left(i_{h-1}\right)}\right)=\left(d_{0}, \ldots, d_{h-1}\right) \notin \rho$ and so $\left(b_{i_{0}}, \ldots, b_{i_{h-1}}\right) \notin \rho^{h^{\kappa}}$. Let $\left(a_{1}, \ldots, a_{h}\right) \in \rho^{h^{\kappa}}$. If $a_{i}=a_{j}$ for some $i \neq j$, then $\left(q\left(a_{1}\right), \ldots, q\left(a_{h}\right)\right) \in \rho$ since $\rho$ is totally reflexive. Otherwise, as $\left(a_{1}, \ldots, a_{h}\right) \in \rho^{h^{\kappa}}$, but $\left(b_{i_{0}}, \ldots, b_{i_{h-1}}\right) \notin \rho^{h^{\kappa}}$ for distinct $i_{0}, \ldots, i_{h-1}$, there is an $1 \leq l \leq h$ such that $a_{l}$ is not equal to any of the $b_{i}$. But then $q\left(a_{l}\right)=u$ by definition of $q$; hence, $\left(q\left(a_{1}\right), \ldots, q\left(a_{h}\right)\right) \in \rho$.
3.4.14 Theorem. If $\rho$ is a central relation, then $\operatorname{Pol}(\rho)$ is a maximal clone.

Proof. Follows from the previous lemma together with Theorems 3.4.1 and 3.4.12.

## $3.5 h$-regularly generated relations

Our next step is to show that $h$-regularly relations generate maximal clones. As those relations are obviously totally reflexive and totally symmetric, the results at the beginning of the last section apply. Our goal is therefore to show the hypotheses of Theorem 3.4.12 are satisfied.
3.5.1 Lemma. Let $D=\left\{d_{1}, \ldots, d_{h}\right\}$ be a subset of $A$ with $\left(d_{1}, \ldots, d_{h}\right) \notin \rho$. Then there is an unary $f \in \operatorname{Pol}(\rho)$ satisfying $f\left(d_{i}\right)=\alpha_{i}, 1 \leq i \leq h$.

Proof. Denote by $\varphi$ the surjection from $A$ onto $h^{\lambda}$ such that $\rho=\varphi^{-1}\left(\omega_{\lambda}\right)$. Set $b_{i}=\varphi\left(d_{i}\right)$ for $1 \leq i \leq h$; then $\left(b_{1}, \ldots, b_{h}\right) \notin \omega_{\lambda}$. That means there is $1 \leq j \leq \lambda$ such that all $b_{i j}$ are distinct, $1 \leq i \leq h$. Thus, the function $s(x)=b_{x j}$ is a bijection from $\{1, \ldots, h\}$ onto $h$. We define $f$ as follows: If $s^{-1}\left(\varphi(x)_{j}\right)=l$, then $f(x)=\alpha_{l}$. Then, as $s^{-1}\left(\varphi\left(d_{i}\right)_{j}\right)=s^{-1}\left(b_{i j}\right)=i$, we get $f\left(d_{i}\right)=\alpha_{i}, 1 \leq i \leq h$. If $\left(a_{1}, \ldots, a_{h}\right) \in \rho$, then $\left(\varphi\left(a_{1}\right), \ldots, \varphi\left(a_{h}\right)\right) \in \omega_{\lambda}$ and so $\varphi\left(a_{i}\right)_{j}=\varphi\left(a_{l}\right)_{j}$ for some $i \neq l$. Thus, $f\left(a_{i}\right)=f\left(a_{l}\right)$ so that $\left(f\left(a_{1}\right), \ldots, f\left(a_{h}\right)\right) \in \rho$ since $\rho$ is totally reflexive. Hence, $f$ preserves $\rho$.

To keep the notation simple, we will sometimes identify the set of tuples $h^{\lambda}$ with its interpretation as a set of natural numbers, sometimes not, whichever is simpler. In our interpretation, we let the numbers $i \in h$ correspond to the tuples $(i, 0, \ldots, 0) \in h^{\lambda}$. Observe that this implies $\left(l_{i}\right)_{1}=l_{i}$ for every $l \in h^{\lambda}$ and every $1 \leq i \leq \lambda$.
3.5.2 Lemma. If we choose any enumeration $\left\{\beta_{0}, \ldots, \beta_{\kappa-1}\right\}$ of $A$ such that $\varphi\left(\beta_{i}\right)=i, i \in h^{\lambda}$, then there is a function $r \in \operatorname{Pol}(\rho)$ which takes all values of $A$ on $\left\{\beta_{0}, \ldots, \beta_{h-1}\right\}$.

Proof. Set $F_{l}=\{x \in A \mid \varphi(x)=l\}$ for $l \in h^{\lambda}$ and denote the elements of $F_{l}$ by $c_{l 0}, \ldots, c_{l n_{l}}$. Set further $n^{*}=\max \left\{n_{l} \mid l \in h^{\lambda}\right\}$ and define for $j \leq n^{*}$ the $n^{*}$-tuple $\operatorname{code}(j)$ to contain $\beta_{1}$ at its $j$-th component and $\beta_{0}$ in all other components. We write $n=\lambda+n^{*}$ and for $i \in h^{\lambda}$ and $j \leq n_{i}$ we define $d_{i j}$ to be the $n$-tuple $\left(\beta_{i_{1}}, \ldots, \beta_{i_{m}}, \operatorname{code}(j)\right)$. The set of all $d_{i j}$ we call $D$. We define an $n$-ary $r$ on $D$ by $r\left(d_{i j}\right)=c_{i j}$, and for $a \in A^{n} \backslash D$ we set $r(a)=\beta_{l}$ where $l=\left(\varphi\left(a_{1}\right)_{1}, \ldots, \varphi\left(a_{\lambda}\right)_{1}\right) \in h^{\lambda} \leq \kappa$. As every element $a \in A$ is for some $l \in h^{\lambda}$ an element of $F_{l}$, we have that $a=c_{l j}$ for some $j \leq n_{l}$ and so $r$ is onto. We claim that for all $a \in A^{n}$ we have that $\varphi(r(a))=\left(\varphi\left(a_{1}\right)_{1}, \ldots, \varphi\left(a_{\lambda}\right)_{1}\right)$. If $a \in A^{n} \backslash D$, then this is a direct consequence of our assumption that $\varphi\left(\beta_{i}\right)=i$ for $i \in h^{\lambda}$. Now if $a=d_{i j} \in D^{n}$ for some $i \in h^{\lambda}$ and $j \leq n_{i}$, then $r(a)=r\left(d_{i j}\right)=c_{i j}$ so that $\varphi\left(r\left(d_{i j}\right)\right)=i$ by the definition of $F_{i}$. On the other hand, $a_{l}=\beta_{i_{l}}, 1 \leq l \leq \lambda$. Thus, again by our assumption on $\varphi, \varphi\left(a_{l}\right)_{1}=\varphi\left(\beta_{i_{l}}\right)_{1}=\left(i_{l}\right)_{1}=i_{l}$. Hence, $\left(\varphi\left(a_{1}\right)_{1}, \ldots, \varphi\left(a_{\lambda}\right)_{1}\right)=\left(i_{1}, \ldots, i_{\lambda}\right)=i=\varphi(r(a))$. We have proven our claim.
We show that $r \in \operatorname{Pol}(\rho)$. Let $\left(r\left(a_{1}\right), \ldots, r\left(a_{h}\right)\right) \notin \rho, a_{i} \in A^{n}, 1 \leq i \leq h$. Then by the definition of a $h$-regularly generated relation, $\left(\varphi\left(r\left(a_{1}\right)\right), \ldots, \varphi\left(r\left(a_{h}\right)\right)\right) \notin \omega_{\lambda}$ which means there is $1 \leq j \leq \lambda$ such that all $\varphi\left(r\left(a_{i}\right)\right)_{j}$ are distinct. By our last claim we have $\varphi\left(r\left(a_{i}\right)\right)_{j}=\varphi\left(a_{i j}\right)_{1}$ and so all $\varphi\left(a_{i j}\right)_{1}$ are distinct, $1 \leq i \leq h$. Hence, by the definition of $\omega_{\lambda},\left(\varphi\left(a_{1 j}\right), \ldots, \varphi\left(a_{h j}\right)\right) \notin \omega_{\lambda}$ and so $\left(a_{1 j}, \ldots, a_{h j}\right) \notin \rho$. Thus, $\left(a_{1}, \ldots, a_{h}\right) \notin \rho^{n}$ and we conclude that $r$ preserves $\rho$.
3.5.3 Lemma. If $D=\left\{d_{1}, \ldots, d_{h}\right\}$ is a subset of $A$ with the property that $\left(d_{1}, \ldots, d_{h}\right) \notin \rho$, then there is an n-ary $q \in \operatorname{Pol}(\rho)$ which takes all values of $A$ on $D^{n}$.

Proof. Let $\left\{\beta_{0}, \ldots, \beta_{\kappa-1}\right\}$ be an enumeration of $A$ with $\varphi\left(\beta_{i}\right)=i, i \in h^{\lambda}$, and let $r \in \operatorname{Pol}(\rho)$ be the function from Lemma 3.5.2. By Lemma 3.5.1 there is a function $g \in \operatorname{Pol}(\rho)$ with $g\left(d_{i}\right)=\beta_{i-1}, 1 \leq i \leq h$. Setting $q\left(x_{1}, \ldots, x_{n}\right)=r\left(g\left(x_{1}\right), \ldots, g\left(x_{n}\right)\right)$ proves the lemma.
3.5.4 Lemma. Let $3 \leq h \leq \kappa$. If $\iota_{h}^{A}$ is $h$-regularly generated, then $\lambda=1$ and $h=\kappa$.

Proof. If $\lambda \geq 2$, then for the vectors $b_{i}=(i, 0, \ldots, 0), 1 \leq i \leq h-1$, and $b_{h}=(1,1,0, \ldots, 0)$ we have $\left(b_{1}, \ldots, b_{h}\right) \in \omega_{\lambda}$. But as those tuples are all distinct, it is impossible that $\iota_{h}^{A}=\varphi^{-1}\left(\omega_{\lambda}\right)$, contradiction. Assume $h<\kappa$. Then, as $\lambda=1, \varphi$ is not one-one and hence there exist distinct $a_{1}, \ldots, a_{h} \in A$ such that $\left(\varphi\left(a_{1}\right), \ldots, \varphi\left(a_{h}\right)\right) \in \omega_{1}$. But this implies $\left(a_{1}, \ldots, a_{h}\right) \in \iota_{h}^{A}$, contradiction.
3.5.5 Lemma. Let $\kappa \geq 3$. If $\rho=\iota_{\kappa}^{A}$, then $\operatorname{Pol}(\rho)$ is a maximal clone.

Proof. $\operatorname{Pol}(\rho)$ contains all unary functions. For if $f \in \mathcal{F}_{1}$ and $\left(a_{1}, \ldots, a_{\kappa}\right) \in A^{\kappa}$ has two identical components, then the same holds for $\left(f\left(a_{1}\right), \ldots, f\left(a_{\kappa}\right)\right)$. Therefore, if $g \notin \operatorname{Pol}(\rho), g$ must depend on at least two variables. But in order to produce a tuple not in $\iota_{\kappa}^{A}, g$ must take all values in $A$. Hence Theorem 3.4.9 yields $\langle\operatorname{Pol}(\rho) \cup\{g\}\rangle=\mathcal{F}$. Observe that $\kappa \geq 3$ is necessary as otherwise $\operatorname{Pol}(\rho)=\mathcal{F}$.
3.5.6 Theorem. Let $3 \leq h \leq \kappa$. If $\rho$ is a h-regularly generated relation, then $\operatorname{Pol}(\rho)$ is $a$ maximal clone.

Proof. If $h=\kappa$, then $\rho=\iota_{\kappa}^{A}$ and $\operatorname{Pol}(\rho)$ is maximal by Lemma 3.5.5. Otherwise $\rho \neq \iota_{h}^{A}$ by Lemma 3.5.4 and application of Lemma 3.5.3 together with Theorem 3.4.12 proves the theorem.

### 3.6 Prime affine relations

Let $\rho$ be a prime affine relation with respect to $(A,+)$. Recall that by definition $(A,+)$ is an abelian group and every $a \in A$ has order $p$, where $p$ is a prime. It is a basic fact from the theory of abelian groups that in this case $|A|=p^{m}$ for some $m>1$. Moreover, $(A,+)$ is isomorphic to the additive group of the field $\mathfrak{G F}\left(p^{m}\right)$ with $p^{m}$ elements. It is for this reason that we can define a multiplication $\cdot$ on $A$ so that $(A,+, \cdot)$ is isomorphic to $\mathfrak{G F}\left(p^{m}\right) .(A,+, \cdot)$ has a primitive element which we call $e$. The neutral elements of + and $\cdot$ we denote by 0 and 1 respectively. In this context, we understand a polynomial to be a function in $\left.<\left\{+, \cdot,(a)_{a \in A}\right\}\right\rangle$. Naturally enough, our approach for proving $\rho$ maximal will be to construct polynomials. We recall the following fact:
3.6.1 Lemma. Every $f \in \mathcal{F}$ is a polynomial. Furthermore, $f\left(x_{1}, \ldots, x_{n}\right)$ can be uniquely expressed as

$$
\begin{equation*}
f\left(x_{1}, \ldots, x_{n}\right)=\sum_{\left(l_{1}, \ldots, l_{n}\right) \in \kappa^{n}} a_{l_{1} \ldots l_{n}} x_{1}^{l_{1}} \ldots x_{n}^{l_{n}} \tag{3.6.1}
\end{equation*}
$$

Proof. It is well-known that every function over a finite field is a polynomial, and it is trivial that $f$ can then be expressed in the form (3.6.1). For the uniqueness, note that there are $\kappa^{\kappa^{n}}$ polynomials of that form which is exactly the number of $n$-ary functions over $A$.
3.6.2 Lemma. The constant functions, the functions $h(x)=a \cdot x, a \in A$, and the operations + and - are affine.

Proof. This is trivial.
3.6.3 Lemma. The functions $h(x)=x^{p^{i}}, 0 \leq i \leq m-1$ are affine.

Proof. We calculate $h(x+y)=(x+y)^{p^{i}}=\sum_{j=0}^{p^{i}}\binom{p^{i}}{j} x^{j} y^{p^{i}-j}$. But $\binom{p^{i}}{j} \equiv 0(p)$ for $1 \leq j \leq p^{i}-1$. Thus, $h(x+y)=x^{p^{i}}+y^{p^{i}}=h(x)+h(y)=h(x)+h(y)-h(0)$ and $h$ is affine.
3.6.4 Corollary. The functions of the form

$$
\begin{equation*}
f\left(x_{1}, \ldots, x_{n}\right)=a_{0}+\sum_{i=1}^{n} \sum_{j=0}^{m-1} a_{i j} x_{i}^{p^{j}} \tag{3.6.2}
\end{equation*}
$$

are affine.
3.6.5 Lemma. If $g \in \mathcal{F}$ is not a function defined by (3.6.2), then $\langle\operatorname{Pol}(\rho) \cup\{g\}\rangle$ contains a function $h(x, y)=\sum_{i, j=0}^{p^{m}-1} a_{i j} x^{i} y^{j}$ with at least one coefficient $a_{s t} \neq 0$, where $1 \leq s, t \leq p^{m}-1$.

Proof. We write $g$ as a polynomial: $g\left(x_{1}, \ldots, x_{n}\right)=\sum_{\left(l_{1}, \ldots, l_{n}\right) \in \kappa^{n}} a_{l_{1} \ldots l_{n}} x_{1}^{l_{1}} \ldots x_{n}^{l_{n}}$. If for one of the coefficients $a_{l_{1} \ldots l_{n}} \neq 0$ there are $1 \leq i, j \leq n, i \neq j$ such that $l_{i}$ and $l_{j}$ are not zero, then setting all variables except $x_{i}$ and $x_{j}$ to 1 yields the desired function. If on the other hand all non-zero coefficients have the form $a_{0 \ldots 0 l_{i} 0 \ldots 0}$, then $g\left(x_{1}, \ldots, x_{n}\right)=g\left(x_{1}, 0, \ldots, 0\right)+$ $g\left(0, x_{2}, 0, \ldots, 0\right)+\ldots+g\left(0, \ldots, 0, x_{n}\right)+c$. Thus there is $1 \leq q \leq n$ with the property that $g\left(0, \ldots, 0, x_{q}, 0, \ldots, 0\right)$ has not the form (3.6.2), for otherwise $g$ would be of that form as well which it is not. Set $f(x)=g\left(0, \ldots, 0, x_{q}, 0, \ldots, 0\right)$ and write $f(x)=\sum_{i=0}^{p^{m}-1} b_{i} x^{i}$. Let $d$ be the greatest index in that sum such that $d$ is not a power of $p$ and $b_{d} \neq 0 ; d=d_{1} \cdot p^{t}$, where $t \geq 0$ and $d_{1} \geq 2$ is not divisible by $p$. Set $h(x, y)=f(x+y)=\sum_{i, j=0}^{p^{m}-1} a_{i j} x^{i} y^{j}$. Then $a_{d-p^{t}, p^{t}}=\binom{d}{p^{t}} b_{d}$. We show that $\binom{d}{p^{t}}$ is not divisible by $p$ :

$$
\binom{d}{p^{t}}=\binom{d_{1} \cdot p^{t}}{p^{t}}=\frac{d_{1} p^{t}\left(d_{1} p^{t}-1\right) \ldots\left(d_{1} p^{t}-p\right) \ldots\left(d_{1} p^{t}-2 p\right) \ldots\left(d_{1} p^{t}-p^{t}+1\right)}{p^{t}\left(p^{t}-1\right) \ldots\left(p^{t}-p\right) \ldots\left(p^{t}-2 p\right) \ldots\left(p^{t}-p^{t}+1\right)}
$$

One readily checks that all factors in the enumerator divisible by powers of $p$ have corresponding factors in the denominator divisible by the same power of $p$. Hence, $a_{d-p^{t}, p^{t}}=\binom{d}{p^{t}} b_{d}$ is not 0 modulo $p$ and the lemma has been proven.
3.6.6 Lemma. If $g \in \mathcal{F}$ is not a function defined by (3.6.2), then $\langle\operatorname{Pol}(\rho) \cup\{g\}>$ contains the function $c(x, y)=x^{s} y^{t}$ for some $1 \leq s, t \leq p^{m}-1$.

Proof. Let $h(x, y)$ be provided by Lemma 3.6.5. If all of the $a_{i j},(i, j) \neq(s, t), 0 \leq i, j, \leq p^{m}-1$, are 0 , we are finished by setting $c(x, y)=a_{s t}^{-1} h(x, y)$. So let $a_{u v} \neq 0$ for $(u, v) \neq(s, t)$, and assume without loss of generality that $u \neq s$. Set $r(x, y)=e^{u} h(x, y)-h(e x, y)$, where $e$ is the primitive element of $(A,+, \cdot)$. Then

$$
r(x, y)=\sum_{i, j=0}^{p^{m}-1}\left(e^{u}-e^{i}\right) a_{i j} x^{i} y^{j}=\sum_{i, j=0}^{p^{m}-1} a_{i j}^{\prime} x^{i} y^{j}
$$

Obviously, $a_{u v}^{\prime}=\left(e^{u}-e^{u}\right) a_{u v}=0$. Furthermore, if $a_{i j}=0$, then also $a_{i j}^{\prime}=0$. On the other hand, as $a_{s t} \neq 0$ and $u \neq s$, we have that $a_{s t}^{\prime}=\left(e^{u}-e^{s}\right) a_{s t} \neq 0$. But this implies that
iteration of this process yields a function $d(x, y)=d_{s t} x^{s} y^{t}$ with $d_{s t} \neq 0$. Hence, we can set $c(x, y)=d_{s t}^{-1} d(x, y)$, and as all operations we used in the process were affine we are finished.
3.6.7 Lemma. Let $g \in \mathcal{F}$ have not the form (3.6.2). Then $\langle\operatorname{Pol}(\rho) \cup\{g\}\rangle=\mathcal{F}$.

Proof. We will show that the function $c(x, y)=x \cdot y$ is an element of $\langle\operatorname{Pol}(\rho) \cup\{g\}>$, for then the definition of the polynomials and Lemma 3.6.2 imply the assertion. By Lemma 3.6.6 $<\operatorname{Pol}(\rho) \cup\{g\}>\operatorname{contains} c(x, y)=x^{s} y^{t}$ with $1 \leq s, t \leq p^{m}-1$. Write $s=h p^{u}$ and $t=l p^{v}$, where $u, v \geq 0$ and $h, l \geq 1$ are not divisible by $p$. By Lemma 3.6.3 the functions $a(x)=x^{p^{m-u}}$ and $b(y)=y^{p^{m-v}}$ are affine; thus, the function $w(x, y)=c\left(x^{p^{m-u}}, y^{p^{m-v}}\right)=\left(x^{p^{m-u}}\right)^{s}\left(y^{p^{m-v}}\right)^{t}=$ $\left(x^{p^{m-u}}\right)^{h p^{u}}\left(y^{p^{m-v}}\right)^{l p^{v}}=\left(x^{p^{m}}\right)^{h}\left(y^{p^{m}}\right)^{l}=x^{h} y^{l}$ is affine as well. Consider $q(x, y)=w(x+1, y+$ 1) $\in<\operatorname{Pol}(\rho) \cup\{g\}>$. We write $q(x, y)=(x+1)^{h}(y+1)^{l}=\sum_{i, j=0}^{p^{m}-1} a_{i j} x^{i} y^{j}$. As $a_{11}=\binom{h}{1}\binom{l}{1}=$ $h l$ and $h, l$ are not divisible by $p$, we conclude that $a_{11} \neq 0$. Now $c(x, y)=x \cdot y \in<\operatorname{Pol}(\rho) \cup\{g\}>$ is an immediate consequence of the proof of Lemma 3.6.6.
3.6.8 Lemma. Let $f \in \mathcal{F}$. Then $f$ is affine if and only if it has the form (3.6.2).

Proof. If $f$ has the form (3.6.2), then it is affine by Lemma 3.6.4. If conversely an $f$ existed which is affine but has not the form (3.6.2), then $\langle\operatorname{Pol}(\rho)\rangle=\langle\operatorname{Pol}(\rho) \cup\{f\}>=\mathcal{F}$ by Lemma 3.6.7, which is absurd.
3.6.9 Theorem. If $\rho$ is a prime affine relation, then $\langle\operatorname{Pol}(\rho)\rangle$ is a maximal clone.

Proof. This is the consequence of Lemmas 3.6.7 and 3.6.8.

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