# Characterization theorems in finite geometry 

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# CHARACTERIZATION THEOREMS IN FINITE GEOMETRY 

H.A. WILBRINK

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## PROEFSCHRIFT

ter verkrijging van de graad van doctor in de technische wetenschappen aan de Technische Hogeschool Eindhoven, op gezag van de rector magnificus, prof.dr. S.T.M. Ackermans, voor een commissie aangewezen door het college van dekanen in het openbaar te verdedigen op dinsdag 1 november 1983 te 16.00 uur

## HENDRIKUS ADRIANUS WILBRINK

geboren te Eindhoven

Dit proefschrift is goedgekeurd door de

PROMOTOR: Prof. dr. J.H. van Lint
CO-PROMOTOR: Prof. dr. J.J. Seidel

## PREFACE

Apart from an introductory chapter this thesis consists of the following five papers.

Nearaffine planes, Geom. Dedicata 12 (1), 53-62.

Finite Minkowski planes, Geom. Dedicata 12 (2), 119-129.

Two-transitive Minkowski planes, Geom. Dedicata 12 (4), 383-395.

A oharacterization of the classical unitals, in: Finite geometries, N.L. Johnson, M.J'. Kallaher \& C.T. Long eds., Marcel Dekker, Lecture notes in pure and applied mathematics 82, New York, 1983.

A characterization of two clases of semi-partial geometries by theix parameters, to appear in simon Stevin.

This last paper was written together with Andries Brouwer. The way we worked together on this paper makes it impossible for me to decide what part of the paper is his and what part is mine.

I would like to express my gratitude to the publishers D. Reidel of Geometriae Dedicata, Marcel Dekker of Finite geometries and J.A. Thas of Simon Stevin for their permission to include these papers in this thesis. I would also like to thank my thesis supervisors Prof. dr. J.H. van Lint and Prof. dr. J.J. Seidel for introducing me to combinatorcs and finite geometry, and for more or less forcing me to write this thesis (I still wonder how they did it). Finally, I have to thank the Mathematical Centre and in particular Andries Brouwer and Arjeh Cohen, for their support and interest in my work during the four fine years I spent there in which period all five papers were written.

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SAMENVATTING

## INTRODUCTION

It is the purpose of this first chapter to introduce the nonexpert mathematician to some of the results and techniques from finite geometry in general, and to each of the five papers which constitute the main part of this thesis, in particular. In each of these five papers a characterization of a finite "incidence structure" is given. Howevex, if one wants to fully understand and appreciate a characterization of any object, it is first necessary to get acquainted with the most basic properties of that object. This is what we shall try to achieve here for the objects discussed in the papers. In addition to this we shall take the opportunity to say something about other theorems characterizing geometries of which ours can be viewed as low dimensional cases.

Basically, characterization theorems in findte geometry fall into four classes. First of all there are the purely geometric characterizations such as the theorem of Veblen and Young characterizing the projective spaces (see section 2), or the Buekenhout-Shult theorem on polar spaces (see section 3). Secondly, there are theorems which use some kind of assumption on the automorphism group of the object in question. The Ostrom-Wagner theorem which we shall discuss in section 1 , is a good example of this. Thirdly, there are the characterizations with the help of a combinatorial property as is the case, for example, in the Dembowski-Wagner theorem which we shall prove in section 2. Finally, it is sometimes possible to characterize geometries if one knows that they are embedded in another geometry (see for example the theorem by Buekenhout-Lefèvre in [6]).

Before we start our discussion a word of warning: the geometries that we shall consider are always assumed to be finite (although for some of the results that we shall state this is really not essential).

## 1. PROJECTIVE PLANES AND AFEINE PLANES

Perhaps the most extensively studled objects in finite geometry are the projective planes. There are several ways to give a definition
of a projective plane. Here we shall adopt one which excludes the degenerate cases and which is easy to generalize to a definition for projective spaces of arbitrary dimension.

DEFINITION. Let $x$ be a set of points and $\mathcal{L}$ a collection of distinguished subsets of $X$ called Zinee. Then $(X, \mathcal{L})$ is called a projective plane if $|\mathcal{L}| \geqq 2$ and the following axioms are satisfied :
(P1) If $x$ and $y$ are distinct points, then there is a unique line $L=x y$ such that $x, y \in L$;
(P2) If $L_{1}$ and $L_{2}$ are distinct lines, then they meet in a unique point; (P3) Every line contains at least 3 points.

The classical models of projective planes are obtained as follows. Let $V$ be a 3-dimensional vector space over $\mathbb{F}_{\mathrm{q}}$, the field of q elements. For $X$ take the set of all 1 -dimensional subspaces of $V$ and for $\mathcal{L}$ the set of all 2-dimensional subspaces of $V$ (more precisely, since we have defined lines to be subsets of $X$, a line is not a 2 -dimensional subspace but the set of all 1-dimensional subspaces contained in a 2-dimensional subspace). It is easy to check that now (P1), (P2) and (P3) are satisfied. Indeed, two distinct 1-spaces span a unique 2-space, two distinct 2-spaces in a 3-space meet nontrivially and a 2 -space over $F_{q}$ contains $\left(q^{2}-1\right) /(q-1)=q+1 \geqq 3$ 1-spaces. The question we are interested in is: are these the only examples of projective planes ? The answer is no. In fact so many different kinds of projective planes are known (see e.g. [8]) that a complete classification seems hopeless. Here we shall content ourselves with one example of a class of projective planes which cannot be obtained from a 3-dimensional vector space. To describe these planes it will be more convenient to work with affine planes. By definition an affine plane is an incidence structure of points and lines satisfying (P1) and
(A2) For every point $x$ and line $L$ such that $x \notin I$, there is exactly one line through $x$ which does not meet $I$;
(A3) There exist three noncollinear points.
It is easy to establish the well-known correspondence between affine planes and projective planes: deleting a line $L_{\infty}$ from a projective plane gives an affine plane and conversely every affine plane can be extended to a
projective plane by adding a line "at infinity". If we follow this procedure for the projective plane associated with the $3-$ dimensional vector space $\left(\mathbb{F}_{q}\right)^{3}$ and with $I_{\infty}$ defined by $z=0$, say, then every point not on $L_{\infty}$ has a unique representation $\langle(x, y, 1)\rangle$ and can therefore be identified, with $(x, y) \in\left(\mathbb{F}_{q}\right)^{2}$. This gives us the familiar affine planes with point set $\left(\mathbb{F}_{q}\right)^{2}$ and with the lines given by an equation $y=a x+b$ or $x=c$. Now it is possible in the above construction to replace the field $\mathbb{F}_{\mathrm{q}}$ by other algebraic structures. For example a quasifield will do as well. Here, a (finite) quasifield is a set $Q$ with two binary operations, + and . say, such that

1) ( $Q,+$ ) is a group with identity 0 ,
2) $(Q \backslash\{0\}, \cdot)$ is a loop with identity 1 ,
3) $x \cdot(y+z)=x \cdot y+x \cdot z$ for all $x, y, z \in Q$,
4) $0 \cdot x=0$ for all $x \in Q$.

It is not hard to show that every quasifield $Q$ yields an affine plane with point set $Q^{2}$ and lines given by an equation $y=a \cdot x+b$ or $x=c$. We shall describe a class of quasifields known as the Andre quasifielda For the set $Q$ take $F_{q} n$ (as a set) and define addition in $Q$ as in $F_{q^{n}}$. Let $A$ be the group of field automorphisms of $F_{q^{n}}$ fixing the subfield $\mathbb{F}_{q}$ of $\boldsymbol{F}_{q^{n}}$ elementwise, and let $N: F_{q}^{*} n^{\rightarrow} F_{q}^{*}$ be the norm map defined by

$$
N(x)=\prod_{\alpha E_{A}} x^{\alpha}, x \in F_{q^{\prime}}^{*}
$$

If $\mu$ is any map from $\mathbb{F}_{q}^{*}$ into $A$ with $\mu(1)=1$, then we can define a multiplication * on $Q$ to make $Q$ into a quasifield as follows

$$
x * y=x y^{\mu(N(x))} \quad(x, y \in Q)
$$

where on the RHS multiplication is in $\mathbb{F}_{q^{n}}$ of course.
We shall now give some of the properties which characterize the projective planes associated with a 3 -dimensional vector space. The first one is probably the best known.

THEOREM 1. A projeotive plane is isomorphic to a projective plane associated with a 3 -dimensional vector space if and only if the following condition holds:
(Desargues' theorem) If $\mathrm{a}_{1}, \mathrm{a}_{2}, \mathrm{a}_{3}$ and $\mathrm{b}_{1}, \mathrm{~b}_{2}, \mathrm{~b}_{3}$ are two triangles such that the tines $a_{1} b_{1}, a_{2} b_{2}$ and $a_{3} b_{3}$ are concurrent, then the points $a_{1} a_{2} n_{1} b_{2}$, $\mathrm{a}_{1} \mathrm{a}_{3} \mathrm{nb}_{1} \mathrm{~b}_{3}$ and $\mathrm{a}_{2} \mathrm{a}_{3} \mathrm{nb}_{2} \mathrm{~b}_{3}$ are collinear (see Figure 1).


Figure 1.

We shall only indicate how Theorem 1 can be proved (for details see e.g. [10] or [16]). The basic idea behind the proof of Theorem 1 is that Desargues' theorem is equivalent to the existence of certain automorphisms of the projective plane (an automorphism of a projective plane is a permutation of the points which induces a permutation of the lines). For example, consider Figure 1 and suppose $\sigma$ is an automorphism fixing x and all the points on $L$. Clearly, since every line through $x$ intersects $L$, all lines through $x$ axe also fixed. If $\sigma$ maps $a_{1}$ to $b_{1}$, then apparently $a_{2}$ is mapped onto $b_{2}$ and $a_{3}$ is mapped onto $b_{3}$; in fact we can determine the image of any point. It is easy to see that Desargues' theorem is equivalent to the existence of this type of automorphisms. Now we have already an algebraic structure associated with our projective plane, namely the group generated by these automorphisms. The special properties of these automorphisms allow us to reconstruct a field $F$ and a 3 -dimensional vector space Vover $F$ from this group in such a way that the projective plane we started with is isomorphic to the projective plane associated with V. Aprojective plane in which Desargues' theorem holds is called a Desarguesian projective plane.

Let us now look at a typical group theoretic characterization of the Desarguesian projective planes.

THEOREM 2. (Ostrom \& Wagner [11]) Let $P=(x, \mathcal{L})$ be a projective plane. If
 projective plane.

Here, 2-transitivity means that for all $x_{1}, x_{2}, y_{1}, Y_{2} \in x, x_{1} \neq x_{2}$, $y_{1} \neq y_{2}$, there is a $y \in \Gamma$ such that $x_{i} Y_{m} y_{i}, i=1,2$. Again we only explain the main ideas of the proof. The trick here is to look at involutions, i.e., automorphisms of order 2. By the 2-transitivity, the even number $|x|(|x|-1)$ divides the order of $\Gamma$ so there exist elements of order 2 in $\Gamma$ (notice that finiteness is really essential here). Let $\sigma$ be an involution. If $x \in X$ and if $x$ is nonfixed, i.e., if $x^{\sigma} \neq x$, then the line $x x^{\sigma}$ is fixed for $\left(x x^{\sigma}\right)^{\sigma}=x^{\sigma} x^{\sigma^{2}}=x^{\sigma} x$. Dually, if $L$ is a nonfixed line, then $L \cap L^{\sigma}$ is a fixed point. From these considerations it follows that the configuration of fixed points and lines of $\sigma$ is either
a) a subplane, or
b) fixes all points on a line $L$ and all lines through a point $x$.

The easy part of the proof is case b), since here $\sigma$ is one of the automorphisms whose existence is equivalent to Desargues' theorem (the only problem here is to show that there are sufficiently many of these automorphisms). The hard part is case b). Suffice it to say that here an induction argument can be used to finish the proof.

We shall see later on that this technique of looking at involutions can also be used to characterize the 2-transitive Minkowski planes.

## 2. PROJECTIVE SPACES

Let $V$ be a vector space of arbitrary dimension. Again we shall use the projective terminology and call the 1 -dimensional subspaces points and the 2-dimensional subspaces lines. Clearly, the points and lines satisfy the axioms (P1) and (P3) of the previous section but (P2) is only satisfied for those lines $L_{1}$ and $L_{2}$ which are contained in a plane (a 3dimensional subspace). In terms of points and lines only, this is expressed
in (P4).
(p4) (Pasch's axiom) If $M_{1}$ and $M_{2}$ are lines meeting in a point $x$ and $L_{1}$ and $L_{2}$ are lines both meeting $M_{1}$ and $M_{2}$ not in $x_{\text {, }}$ then $\mathrm{L}_{1}$ and $\mathrm{L}_{2}$ meet.


DEFINITION. Let $X$ be a set of points and $I$ a collection of distinguished subsets of X called lines. Then ( $\mathrm{X}, \mathrm{S}$ ) is called a projective space if ( P 1 ), (P3) and (P4) are satisfied.
clearly, every projective plane is a projective space. The following theorem, due to Veblen \& Young, shows that for higher dimensions there is no analogue to the "nondesarguesian" planes.

THFOREM 3. Let ( $\mathrm{X}, \mathrm{L}$ ) be a projective space containing two nonintersecting Zines. Then ( $\mathrm{X}, \mathcal{L}$ ) is isomorphic to the geometry of 1 - and 2 -dimensional subspaces of a vector space.

We explain the main steps in the proof of this theorem. Let ( $\mathrm{X}, \mathrm{L}$ ) be a projective space. A subset $Y \subset X$ is called a subspace if every line which meets $Y$ in at least two points, is completely contained in $Y$. Clearly every subspace together with the lines it contains is also a projective space. It is also easy to prove that if $Y$ is any subspace and $x$ is any point not contained in $Y$, then the set $Z$ of all points on lines through $x$ which meet $Y$ (i.e., $Z=Y \underset{Y}{ } \quad X y$ ) is also a subspace. If we take for $Y$ a line, the resulting $Z$ is easily seen to be a projective plane. Now look at Figure 1, not as a configuration in the plane but with $x$ not in the plane generated by $a_{1}, a_{2}$ and $a_{3}$, say. The points $a_{i} a_{j} \cap b_{i} b_{j}, 1 \leqq i<j \leq 3$, are all on the intersection line $L$ of the planes generated $b y a_{1}, a_{2}, a_{3}$ and $b_{1}, b_{2}$, $b_{3}$, so Desargues' theorem holds in this case. In fact, Desargues' theorem holds in all cases for, if $x, a_{1}, a_{2}, a_{3}$ happen to be in a plane, we can always view the configuration as the projection of a nonplanar configuration from a point onto the plane generated by $x, a_{1}, a_{2}, a_{3}$. By Theorem 1 we now know already that all projective planes which are properly contained in a projective space are isomorphic to a projective plane associated with a

3-đimensional vector space (this result is an example of a characterization using an embeddability property). The rest of the proof consists in glueing together these 3-dimensional vector spaces to one big vector space (see e.g. [3],[12] or [16]for more details).

As a typical application of Theorem 3 we shall prove the DembowskiWagner theorem which is a combinatorial characterization of projective spaces in terms of points and hyperplanes. For this we need some terminology which will also be useful later on. A t-design with parameters $v_{r}$ $k, \lambda$ (or a $t-(v, k, \lambda)$ design) is a pair ( $\mathrm{X}, \mathrm{B}$ ) where $B$ is a collection of $k$-subsets (called blocks) of a set $X$ of $v$ points such that every t-subset of $X$ is contained in exactly $\lambda$ blocks. For any two points $x$ and $y$ in $a$ 2-design we define the line through $x$ and $y$ as the intersection of the blocks containing $x$ and $y$. Notice that every two distinct points in a 2-design are on a unique line. For example, let $v$ be a vector space of dimension $n$ over ${ }^{I}{ }_{q}$, $X$ the set of all points of the projective space associated with $V$ and let $B$ be the set of all hyperplanes of $V$. Then ( $X, B$ ) is a $2-\left(\frac{q^{n}-1}{q-1}, \frac{q^{n-1}-1}{q-1}, \frac{q^{n-2}-1}{q-1}\right)$ design and the lines in the 2 -design sense are precisely the lines in the projective space sense. This design has the property that the total number of blocks is equal to the total number of points. A 2-design with this property is called symmetric or projective.

THEOREM 4. (Dembowski-Wagnex) Let ( $\mathrm{X}, \mathrm{B}$ ) be a symmetric 2-(v,k, $\lambda$ )design. Then $(\mathrm{X}, \mathrm{B})$ is the design of points and hyperplanes of a projective space if and only if every line has at least $(\mathrm{v}-\lambda) /(\mathrm{k}-\lambda)$ points.

PROOF. Since $|B|=v$, every point is on $k$ blocks. Let $L$ be any line. Since $L$ is contained in $\lambda$ blocks, every point $x$ on $L$ is on $k-\lambda$ blocks $B$ such that L $\cap B=\{x\}$.Therefore $v-(\lambda+|L|(k-\lambda))$ blocks do not meet $L$. From our hypothesis it follows that $|L|=(v-\lambda) /(k-\lambda)$ and that every line meets every block. Let $x$ be any point not on $L$ and suppose that $\rho$ blocks contain $L$ and $x$. Then $k-\rho$ blocks contain $x$ but not $L$. This number also equals $|L|(\lambda-\rho)$ (for each $y \in L$ there are $\lambda-\rho$ blocks $B$ on $x$ and $y$ such that $L \cap B=\{y\})$. Therefore $k-\rho=|L|(\lambda-\rho)$ and so $\rho$ is a constant. Define planes as the intersection of all blocks containing three noncollinear
points. Any three noncollinear points now determine a unique plane. Let L and $M$ be two distinct lines in a plane E. Let $B$ be a block containing $L$ but not $M$. Then $L=B \cap E$, so $L \cap M=(B \cap E) \cap M=B \cap(E \cap M)=B \cap M \neq \phi$, i.e. any two lines in a plane meet. This proves Pasch's axiom.

## 3. SYMPLECTIC, UNITARY AND ORTHOGONAL GEOMETRY

We shall now turn to certain substructures of projective spaces for which there is a characterization quite similar to the characterization of Veblen \& Young for projective spaces. Let us start with an analytic description of these substructures. Suppose $V$ is a vector space of dimension $n$ over $\mathbb{F}_{q}$ and let $\sigma$ be an automorphism of $\boldsymbol{F}_{q}$. We shall often write $\bar{\lambda}=\lambda^{0}$ for $\lambda \in \mathbb{F}_{\mathrm{q}}$. A ( $\sigma$-sesquilinear ) form $f$ on $V$ is a map $\mathrm{f}: \mathrm{V} \times \mathrm{V} \rightarrow \mathrm{F}_{\mathrm{q}}$ satifying
i) $f(\lambda x, y)=\lambda f(x, y)$ and $f(x, \lambda y)=\lambda f(x, y), x, y \in v, \lambda \in \mathbb{F}_{q}$;
ii) $f(x, y+z)=f(x, y)+f(x, z)$ and $f(x+y, z)=f(x, z)+f(y, z), x, y, z \in V$.

The form $f$ is called reflexive if for all $x, y \in V, f(x, y)=0 \Leftrightarrow f(y, x)=0$ and $f$ is called nondegenerate if $f(x, y)=0$ for all $x \in V \Rightarrow y=0$. If $n \geqq 2$ and $f$ is a nondegenerate reflexive form on $V$, then there are only a few possibilities for $f$ (see e.g. [2]) :
i) $\sigma=1$ and $f(x, x)=0$ for all $x \in v$.

In this case fiscalled a symplectic form and it is possible to show that $n$ has to be even and that w.r.t. to a suitable basis $v_{1}, v_{2}, \ldots, v_{n}$ of $V$,
$f(x, y)=\xi_{1} n_{2}-\xi_{2} \eta_{1}+\xi_{3} \eta_{4}-\xi_{4} \eta_{3}+\cdots+\xi_{n-1} \eta_{n}-\xi_{n} \eta_{n-1}, x=\sum \xi_{i} v_{i}, y=\sum \eta_{i} v_{i}$.
ii) $\sigma^{2}=1, \sigma \neq 1$ and for some $\lambda_{0} \in \mathbb{F}_{q}, \lambda_{0} f(x, y)=\overline{\lambda_{0} f(y, x)}$ for all $x, y \in v$. In this case $\lambda_{0} f$ is called hermitian and w.r.t. a suitable basis
$v_{1}, v_{2}, \cdots, v_{n}$ of $v$

$$
f(x, y)=\Sigma \xi_{i} \eta_{i}, x=\Sigma \xi_{i} v_{i}, y=\Sigma \eta_{i} v_{i}
$$

iii) $\sigma=1$ and $f(x, y)=f(y, x)$ for all $x, y \in V$.

In this case fis called symmetric. For even $q$, symmetric forms are not very interesting and for odd $q$, symmetric forms are equivalent with quadratic forms which we shall now discuss.

A quadratic form $Q$ on $V$ is a map $Q: V \rightarrow F_{q}$ such that
a) $Q(\lambda x)=\lambda^{2} Q(x)$ for all $\lambda \in \mathbb{F}_{q^{\prime}} x \in V$, and
b) $f(x, y):=Q(x+y)-Q(x)-Q(y)$ defines a bilinear form on $v$.

Notice that $f$ is symmetric and that $f(x, x)=Q(2 x)-2 Q(x)=2 Q(x)$. Conversely if $q$ is odd and $f$ is any symmetric form on $V$, then $Q(x):=\frac{1}{2} f(x, x)$ is a quadratic form with associated bilinear form $f$, so for $q$ odd, $f$ and $Q$ determine each other. A quadratic form $Q$ is called nondegenerate if $Q(x) \neq 0$ for all $x \in \mathbb{X}\{0\}$ which satisfy $f(x, y)=0$ for all $y \in V$ (for odd $q$ this is equivalent to $f$ is nondegenerate, but if $q$ is even $f$ can be degenerate whereas $Q$ is not (see type (I) below)). The standard forms for a nondegenerate quadratic form w.r.t. a suitable basis are as follows. If $n$ is odd there is essentially one type:
(I) $Q(x)=\xi_{1} \xi_{2}+\xi_{3} \xi_{4}+\cdots+\xi_{n-2} \xi_{n-1}+\alpha \xi_{n}^{2}$, for some $\alpha \in F_{q}$.

If n is even there are two types:
(II) $Q(x)=\xi_{1} \xi_{2}+\xi_{3} \xi_{4}+\cdots+\xi_{n-1} \xi_{n}$, or
(III) $Q(x)=\xi_{1} \xi_{2}+\xi_{3} \xi_{4}+\cdots+\xi_{\mathrm{n}-3} \xi_{\mathrm{n}-2}+\xi_{\mathrm{n}-1}^{2}+\alpha \xi_{\mathrm{n}-1} \xi_{\mathrm{n}}+\beta \xi_{\mathrm{n}}{ }^{2}$
where $X^{2}+\alpha x+\beta$ is irreducible over $\mathbb{F}_{q}$.
Suppose $f$ is a reflexive form on V. If $f(x, y)=0$ we write $x \perp y$ and say that $x$ and $y$ are orthogonal. Since $£$ is reflexive, $\perp$ is a symmetric relation. For $X \subset V$ we set

$$
x^{\perp}:=\{v \in v \mid v \perp x \text { for all } x \in X\}
$$

$A$ subspace $x$ of $V$ is called totally $i s o t r o p i c$ if $X \subset X^{\perp}$, i.e. if $f(x, y)=0$ for all $x, y \in X$. Similarly, if $Q$ is a quadratic form on $V$, then any subspace $X$ with $Q(X)=0$ is called totally singutar. (If $q$ is odd, then $X$ is totally singular if and only if $X$ is totally isotropic w.r.t. the bilinear form $f$ associated with $Q$.$) A vector space V$ equipped with a nondegenerate symplectic, hermitian or quadratic form is called a symplectic,unitary or orthogonal geometry. Especially the set of all totally isotropic (singular) points in symplectic, unitary and orthogonal geometry gives us all kinds of interesting configurations. for example, take a quadratic form of type (III) with $n=4$ and work over $\mathbb{R}$ for the moment with $Q(x)=\xi_{1} \xi_{2}+\xi_{3}^{2}+\xi_{4}^{2}$.

The set of totally singular points here is a sphere (put $\xi_{1}=\eta_{1}+\eta_{2}$, $\xi_{2}=\eta_{1}-\eta_{2}$ and look in the affine 3 -space defined by $\eta_{2}=1$, so any three totally singular points determine a plane which will intersect the sphere in a conic. precisely the same is true over a finite field: let $x$ be the set of totally singular points and $B=\{X \cap E \mid E$ a plane with $|X \cap E| \geqq 3\}$, then ( $X, B$ ) is a 3-design. Keeping the picture of the sphere in mind it is easy to compute the parameters of the design. If $p$ is any totally singular point, then $P$ is on $q+1$ tangent lines (all the lines in the plane tangent to the sphere passing through p) which carry no other points of the sphere, and therefore on $\left(q^{2}+q+1\right)-(q+1)=q^{2}$ lines which intersect the sphere in one other point. Hence $|x|=q^{2}+1$, and a similar argument in the plane shows that every conic contains $q+1$ points. Thus ( $X, B$ ) is a $3-\left(q^{2}+1, q+1,1\right)$ design. A Möbius plane is by definition a $3-\left(n^{2}+1, n+1,1\right)$ design. The Möbius planes that we have just constructed are characterized by the fact that they satisfy the Theorem of Miquel (see [18]). They play a role similar to that of the Desarguesian planes in the theory of projective planes. Here also, "nonmiquelian" Mobius planes are known to exist (although not as many as nondesarguesian projective planes). A similar story can be told by starting off with a quadratic form $Q(x)=\xi_{1} \xi_{2}+\xi_{3} \xi_{4}$ of type (II) . We then arrive at the so-called Minkowski planes which we shall discuss in greater detail in the next section.

There is a very satisfactory characterization of the symplectic, unitary and orthogonal geometries which have totally isotropic or totally singular subspaces of dimension at least three, known as the BuekenhoutShult theorem, which we shall now formulate.

DEFINITION. Let $X$ be a set of points and $\mathcal{L}$ a collection of distinguished subsets of X'called lines such that
i) the set of lines is nonempty and each line has at least three points,
ii) no point is collinear with all remaining points,
iii) for every point $x$ and every line $L$ not containing $x, x$ is collinear with either one or all points of $L$.

Then ( $\mathrm{X}, \mathcal{L}$ ) is called a polar space.

Every symplectic, unitary or orthogonal geometry containing totally isotropic (totally singular) lines yields a polar space in the following way: points are the totally isotropic (singular) points, lines are the totally isotropic (singular) lines. Let us check iii) for a symplectic or unitary space $v$. Let $\langle x\rangle$ be a totally isotropic point and Latotally isotropic line. Since $\langle x\rangle^{\perp}=\{y \mid f(x, y)=0\}$ is a hyperplane of $V$, the 2 -dimensional subspace $L$ intersects $\langle x\rangle \perp$ nontrivially. If $y \in L \cap\langle x\rangle \perp, y \neq 0$, then $f(\lambda x+\mu y, o x+O y)=0$ since $f(x, x)=f(y, y)=f(x, y)=0$, so the line $\langle x, y\rangle$ is totally isotropic. If $L \notin\langle x\rangle \perp$, then $\langle x\rangle$ is collinear (in the polar space sense) with exactly one point of $L$, if $L \subset\langle x\rangle \perp$, then $\langle x\rangle$ is collinear with all points of $L$.

THEOREM 5. Let $(\mathrm{X}, \mathrm{L})$ be a polar space. Then
a) ( $\mathrm{X}, \mathrm{L}$ ) is isomorphic to the geometry of all totally isotropic or totally singular points and lines of a symplectic, unitary or orthogonat
geometry, or
b) ( $\mathrm{x}, \mathrm{L}$ ) satisfies the following stronger version of $i t i$ ):
iv) for every point $x$ and every line $L$ not containing $x, x$ is collinear with excotly one point of L .

The first characterization of polar spaces was obtained by Veldkamp [17] who used a more complicated set of axioms. This set of axioms was later simplified by Tits (see [15]) and Buekenhout and Shult (see [5]). A polar space which satisfies iv) is called a generalized quadrangle. Here the generalized quadrangles play a role similar to that of the projective planes in the theory of projective spaces. Again many generalized quadrangles are known which are not isomorphic to the geometry of totally isotropic or totally singular points and lines of a symplectic, unitary or orthogonal geometry. For example, the following geometry of ab points and $a+b$ lines as shown in Figure 2 is a generalized quadrangle.


Figure 2.

However，since lines in a projective space over $F_{q}$ carry $q+1$ points， this can only be a geometry of totally isotropic or totally singular points and lines if $a=b=q+1$ for some prime power $q$ ．The orthogonal geometry over $\mathbb{F}_{q}$ of type（II）for $n=4$ belonging to the quadratic form $\mathrm{Q}(\mathrm{x})=\xi_{1} \xi_{2}+\xi_{3} \xi_{4}$ yieldsa generalized quadrangle of this type with $a=b=$ $=q+1$ ；the two sets of $q+1$ mutually disjoint lines correspond to the two sets of rulings on the hyperboloid $\xi_{1} \xi_{2}+\xi_{3} \xi_{4}=0$ ．Additional axioms are necessary to characterize the classical generalized quadrangles．For example，there is a theorem by Buekenhout \＆Lefèvre（see［6］）which says that a generalized quadrangle which is embedded in a projective space is classical．Characterizations using certain（transitivity）properties of the automorphism group have been given by Tits［14］and Walker［19］．Thas and payne（see e．g．［13］）have given a number of characterizations based on geometic and combinatorial assumptions．

## 4．SUMMARY OF THE FIVE PAPERS

The first paper［A］is on nearaffine planes．Nearaffine planes（and more generally nearaffine spaces）were introduced by J．Andre（see e．g．［1］） to describe geometrically vector spaces over nearfields．By definition a nearfietd（ $F,+,+$ ）is a quasifield（as defined in section 1 ）with the additional property that（ $F \backslash\{0\}, \cdot$ ）is a group．Let（ $F,+, \cdot$ ）be a nearfield and set $V=F^{2}$ ．With addition and scalar multiplication on the left（by elements of $F$ ）defined componentwise on $V, V$ is called a vector space of dimension 2 over $F$ ．For $x, y \in V, x \neq y$ ，define the line $x H y$ from $x$ to $y$ by

$$
x 山 y:=F \cdot(y-x)+x
$$

If $F$ happens to be a field，then $V$ is just the standard 2－dimensional vector space over $F$ and the lines $x 山 y$ coincide with the ordinary lines in the Desarguesian affine plane．If $F$ is a proper nearfield，then in general $u, v \in x 山 y$ does not imply $x 山 y=u \| v$ and a rather complicated set of axioms is necessary to describe this geometry．The axioms for a near－ affine plane are chosen in such a way that we get the ordinary affine planes back if the additional property $x L y=y 山 x$ holds for all $x, y \in V$ ．

What we do in this paper is to set up a theory for nearaffine planes which generalizes the theory of translation planes, i.e. affine planes which can be coordinatized by a quasifield in the sense as described in section 1. This leads us to what we have called nearaffine tranelation planee. As for ordinary translation planes, it is possible to give equivalent algebraic, geometric and group theoretic descriptions of nearaffine translation planes. For us nearaffine planes are especially important due to certain connections with Minkowski planes, the subject of papers [B] and [C] which we shall now discuss.

Consider the hyperboloid in projective 3 -space over $\mathbb{F}_{q}$, i.e. the set of totally singular points of the quadratic form $Q(x)=\xi_{1} \xi_{2}+\xi_{3} \xi_{4}$ on $F_{q}^{4}$. The picture to keep in mind here is that of the hyperboloid $x^{2}-y^{2}+z^{2}=1$ (use the transformation $\xi_{1}=x-y, \xi_{2}=x+y, \xi_{3}=z-t, \xi_{4}=z+t$ and take $t=1$ ). There are two families $\mathcal{L}^{+}$and $\mathcal{L}^{-}$of totally singular lines on the hyperboloid. Explicitly these lines are (in $\xi$-coordinates)

$$
\begin{aligned}
& e_{a, b}^{+}:=\langle(a, 0, b, 0),(0, b, 0,-a)\rangle \text { and } \\
& e_{a, b}^{-}:=\langle(a, 0,0, b),(0, b,-a, 0)\rangle
\end{aligned}
$$

where $a, b \in F_{q}$ and at least one of $a$ and $b$ is not equal to zero. We have already pointed out that the totally singular lines form the rather trivial structure of $a(q+1) \times(q+1)$ grid (see fig.2). To obtain an interesting geometry we proceed as in the case of the MObius planes and add the conic intersections of the planes with the set of totally singular points as objects to our geometry. These plane sections are called circles. Any three distinct points on the hyperboloid with the property that no two are on a totally singular line determine a unique plane and therefore a unique circle. In this way we arrive at an icidence structure with a set $M$ of points, two collections $\mathcal{L}^{+}$and $\mathcal{L}^{-}$of subsets of $M$ called innes, and a collection $C$ of subsets of $M$ called circles satisfying the following axioms.
(M1) $\mathcal{L}^{+}$and $\mathcal{L}^{-}$are partitions of $M_{\text {, }}$
(M2) $\left|\ell^{+} \cap \ell^{-}\right|=1$ for all $\ell^{+} \in \mathcal{L}^{+}, \ell^{-} \in \mathcal{L}^{-}$,
(M3) any three points, no two on a line, determine a unique circle $c \in C$,
(M4) $|\ell \cap c|=1$ for all $\ell \in \mathcal{L}^{+} \cup \mathcal{L}^{-}, c \in C$,
(M5) there exist three points, no two of which are on a line.

Such an incidence structure is called a Minkowski plane. Let us prove some elementary properties of Minkowski planes. From (M1) and (M2) it follows that $\left|\ell^{+}\right|=\left|\mathcal{L}^{-}\right|$for all $\ell^{+} \in \mathcal{L}^{+}$and $\left|\ell^{-}\right|=\left|\mathcal{L}^{+}\right|$for all $\ell^{-} \in \mathcal{L}^{-}$. By (M1) and (M4) we have $\left|\mathcal{L}^{+}\right|=|c|$ and $\left|\mathcal{L}^{-}\right|=|c|$ for all $c \in C$. Since $C \neq \varnothing$ by (M3) and (M5), we have proved that $\left|\mathcal{L}^{+}\right|=\left|\mathcal{L}^{-}\right|=|\ell|=|c|$ for all $\ell \in \mathcal{L}^{+} \cup \mathcal{L}^{-}, c \in C$. The number $n:=|c|-1$ is called the order of the Minkowski plane. It is often convenient to think of the points and lines of a Minkowski plane as being arxanged in an $(n+1) \times(n+1)$ square grid.


Figure 3.

Every circle then corresponds to a transversal of this grid intersecting each horizontal and vertical line exactly once. An important property (which for infinite Minkowski planes is an additional axiom) is
(M6) given a circle $c$, a point $P \in c$ and a point $Q \& c, P$ and $Q$ not on $a$ line, there is a unique circle $d$ such that $P, Q \in d$ and $c \cap d=\{p\}$.

To prove this, note that the two noncollinear points P and $Q$ are on $n-1$ circles (Figure 3 shows that there are ( $n-1)^{2}$ points not collinear with $P$ or $Q$ : each circle through $p$ and $Q$ contains $n-1$ of these). Since there are $n-2$ points on $c$ not equal to $p$ and noncollinear with $Q$, there must be exactly $(n-1)-(n-2)=1$ circle through $P$ and $Q$ which does not intersect $c$ in a point distinct from $p$. With the help of (M6) it is not very hard to see that with every point $z$ of a Minkowski plane we can associate an affine plane (the derived plane at $Z$ ) as follows. The points of the affine plane
are the points which are not collinear with $Z$. The lines of the affine plane are the lines of the Minkowski plane missing $Z$ and the circles containing Z . Axiom (A2) for affine planes now corresponds to (M6). In the hyperboloid model this affine plane is clearly visible if we use stereographic projection from $Z$ onto a plane.

It is possible to construct Minkowski planes which are not isomorphic to a Minkowski plane associated with a quadratic form on $\mathbb{T}_{\mathbf{q}}{ }^{4}$. In $[B]$ we show that the known Minkowski planes are characterized by the fact that a certain geometrical condition (called (D) in [B]) holds. The idea behind the proof of this lies in the observation that with any point $z$ of one of the known Minkowski planes we can also associate a nearaffine plane. The points of the nearaffine plane are again the points which are not collinear with $Z$. The lines of the nearaffine plane correspond to the lines and circles missing $Z$. Viewed in this way, condition (D) is nothing but a special case of Desargues' theorem in the nearaffine plane. One can show that (D) implies that all nearaffine planes are nearaffine translation planes. The automorphisms of the nearaffine planes extend to automorphisms of the Minkowski plane. These in turn enable one to reconstruct the algebraic representation of the known Minkowski planes.

In [c] we have generalized the theorem of Ostrom \& Wagner for projective planes (Theorem 2) and Hering's result for Möbius planes (see [9]) to Minkowski planes: if the automorphism group of a Minkowski plane is transitive on pairs of noncollinear points, then the plane is one of the known Minkowski planes. The technique used here is very much the same as in the proof of the Ostrom \& Wagner theorem. Again the basic tool is to study involutions in the automorphism group. Here some rather deep group theory is necessary to reduce to the case where there is an involution which has a subplane as a set of fixed points. Once this is achieved, induction is possible to finish the proof.

In [D] we have characterized the unitary geometry on $F_{q^{2}}^{3}$ which we shall now describe in some detail. Let $q$ be a prime power and $V=\mathbb{F}_{q^{2}}{ }^{*}$ Define a nondegenerate hermitian form ( , ) on $V$ by

$$
(x, y)=\xi_{1} \eta_{1}+\xi_{2} \eta_{2}+\xi_{3} \eta_{3},
$$

for $x=\left(\xi_{1}, \xi_{2}, \xi_{3}\right), y=\left(\eta_{1}, \eta_{2}, \eta_{3}\right) \in v$. Here $T=\lambda^{q}$ for all $\lambda \in \mathbb{F}_{q^{2}}$.

Let $U$ be the set of totally isotropic points, i.e.

$$
u=\{\langle x\rangle \mid(x, x)=0, x \in V \backslash\{0\}\}
$$

Let $\langle x\rangle \in U$ and let $\langle y\rangle$ be any any other point. A point $\langle\lambda x+y\rangle$ on the line $\langle x, y\rangle$ joining $\langle x\rangle$ and $\langle y\rangle$ is in $u$ if $0=(\lambda x+y, \lambda x+y)=\operatorname{Tr}(\lambda(x, y))+(y, y)$, where $\operatorname{Tr}: \mathbb{F}_{q^{2}}+\mathbb{F}_{q}$ is the trace map given by $\operatorname{Tr}(\alpha)=\alpha+\bar{\alpha}, \alpha \in \mathbb{F}_{q^{2}}$. We claim that it is impossible that all points $\langle\lambda x+y\rangle$ are in $u$, i.e., that $(x, y)=0$ and $(y, y)=0$. Suppose on the contrary that $(x, y)=(y, y)=0$. Take any point $\langle z\rangle$ not on the line $\langle x\rangle^{\perp}:=\left\{x^{\prime} \mid\left(x, x^{\prime}\right)=0\right\}$. Then

$$
u=-\frac{(y, z)}{(x, z)}+y
$$

satisfies $(u, x)=(u, y)=(u, z)=0, s o(u, v)=0$ for all $v \in V$, a contradiction. This shows that $\langle x\rangle$ is the only point of $u$ on the line $\langle x\rangle{ }^{\perp}$ and that every other line through $\langle x\rangle$ contains $q$ points $\neq\langle x\rangle$ of $U$ (for $T r$ is an $\mathbb{F}_{q}$-linear map with a kernel of dimension 1 , so $\operatorname{Tr}(\lambda(x, y))=-(y, y)$ has $q$ solutions $\lambda$ if $(x, y) \neq 0)$. Since there are $q^{2}+1$ lines through $\langle x\rangle$, one of which is $\langle x\rangle{ }^{1}$, it follows that $|\mathrm{U}|=1+q \cdot q^{2}=1+q^{3}$. Also, every two distinct points of $U$ are on a unique line of $q+1$ U-points, i.e., we have constructed a $2-\left(q^{3}+1, q+1,1\right)$ design. A $2-\left(n^{3}+1, n+1,1\right)$ design is called a unital ( $n \in \mathbb{N}$ ). For $q=2$ the $2-(9,3,1)$ design is the unique affine plane of order 3. But already for $q=3$ numerous $2-(28,4,1)$ designs are known (see Brouwer [4]) and so we are left with the question what properties are characteristic for the unitals associated with a unitary geometry. It is conjectured that the following "anti-Pasch" axiom will do:

No four distinct points intersect in six distinct lines.
It is easy to show that this property holds for the classical unitals. Suppose $\langle x\rangle,\langle y\rangle,\langle a\rangle,\langle b\rangle,\langle c\rangle,\langle d\rangle$ are six distinct points of $u$ such that they form the configuration of Figure 4.


Figure 4.

Since $a, b, c$ and $d$ are linearly dependent, we may assume that

$$
a+b+c+d=0
$$

and therefore also that $x=a+c, y=a+b$. From $(x, x)=0$ it follows that $(a, c)+(c, a)=0$. Similarly, $(a, b)+(b, a)=0 \quad($ from $(y, y)=0)$ and $(b, c)+(c, b)=0$ (from $(d, a)=0$ and the other relations). Since $a, b$ and $c$ are linearly independent the Gram matrix

$$
\left(\begin{array}{ccc}
0 & (a, b) & (a, c) \\
(b, a) & 0 & (b, c) \\
(c, a) & (c, b) & 0
\end{array}\right)
$$

is nonsingular. Hence $0 \neq(a, b)(b, c)(c, a)+(a, c)(b, a)(c, b)$. This contradicts the other relations.

In [D] we have characterized the classical unitals under additional geometric assumptions. The basic steps in the proof are as follows. Using nontrivial group theory it is easy to prove that once the automorphism group of the unital is large enough, we can only have a classical unital. The geometrical conditions we impose ensure the existence of such an automorphism group. More precisely, for the classical unital we have for $\langle x\rangle \in U$ that the linear transformation

$$
v \mapsto v+\alpha(x, v) x, v \in v
$$

respects the hermitian form (, ) if $\operatorname{Tr}(\alpha)=0$ and so acts as an automorphism of the unital fixing all lines through <x>. These transformations are called the unitary transvections. The geometrical conditions imply the existence of all possible unitary transvections and these generate a 2-transitive group of automorphisms.

We conclude with a discussion of the last paper [E] on semi-partial geometries. The concept of generalized quadrangle has been generalized in a number of ways by replacing the key axiom iv) as formulated in Theorem 5, by a similar axiom. Most of these axioms can be formulated as:

For every point $x$ and every line $E$ with $x \in L$,

$$
\mid\{y \in L \mid x \text { and } y \text { collinear }\} \mid \in s
$$

where $S$ is some finite subset of $N U\{0\}$. By taking $S=\{0, \alpha\}$ one gets the
essential axiom for a semi-partial geometry (for a complete definition see [E]). In this paper we show that certain semi-partial geometries are already determined by some numerical data. There are two cases to consider, namely $\mu=\alpha^{2}$ and $\mu=\alpha(\alpha+1)$ in the notation of [E]. The line of proof in both cases is essentially identical and roughly reads as follows. By results of Debroey [7] it suffices to show that the points and lines of such a semi-partial geometry satisfy the dual of the axiom of pasch (for obvious reasons called the diagonal axiom) . For both the conditions $\mu=\alpha^{2}$ and $\mu=\alpha(\alpha+1)$ there is a straightforward geometric interpretation. The hard part of the proof consists in using this over and over again to show that any two intersecting lines generate a well-behaved "subspace". once this has been achieved it is no longer hard to show that the diagonal axiom holds provided the semi-partial geometry properly contains such a subspace.

## REFERENCES

1. ANDRE, J., On finite non-commutative affine spaces, in: Combinatorics (Part 1), M. Hall \& J.H. van Lint eds., Math. Centre Tracts 55 (1974), 60-107.
2. ARTIN, E., Geometric algbra, Interscience, New York, 1957.
3. BAER, R., Linear algebra and projective geometry, Academic Press, New York, 1952.
4. BROUWER, A.E., Some unitals on 28 points and their embeddings in projective planes of order 9, Math. Centre Report 155/81 (1981)
5. BUEKENHOUT, F. \& E. SHULT, On the foundation of polar geometry, Geom. Dedicata 3 (1974), 155-170.
6. BUEKENHOUT, F. \& C. LEFEVRE, Generalized quadrangZes in projective spaces, Arch. Math. 25 (1974), 540-552.
7. DEBROEY, I., Semi partial geometries satifying the diagonal axiom, J. Geometry 13 (1979), 171-190.
8. DEMBOWSKI, P., Finite geometries, Springer-Verlag, New-York, 1968.
9. HERING, C., Endliche zweifach transitive Möbiusebenen ungerader ordnung, Arch. Math. 18 (1967), 212-216.
10. HUGHES, D. \& F. PIPER, Projective planes, Springer-Verlag, New York, 1973.
11. OSTROM, T.G. \& A. WAGNER, On projective and affine planes with transitive collineation groups, Math. Z. 71 (1959), 186-199.
12. TAMASCHKE, O., Projective Geometrie I \& II, B.I., Mannheim, 1969.
13. THAS, J.A. \& S.E. PAYNE, Classical finite generalized quadrangles : a combinatorial study, Ars Combin., 2 (1976), 57-110.
14. TITS, J., Classification of buildings of spherical type and Moufang polygons : a survey, Atti Coll. Geom. Comb. Roma (1973).
15. TITS, J., Buildings of Spherical Type and Finite BN-Pairs, SpringerVerlag, Lectures Notes \#386, Berlin, 1974.
16. VEBLEN, O. \& J.W. Young, Projective geometry I \& II, Ginn, Boston, 1916.
17. VELDKAMP, F.D. Potar geometry I-V, Proc. of the KNAW (A) 62 (1959), $512-551 ; 63(1960), 207-212$.
18. VAN DER WAERDEN, B.L. \& L.J. SMID, Eine Axiomatik der Kreisgeometrie und der Laguerre-Geometrie, Math. Ann. 110 (1935), 753-776.
19. WALKER, M., On the stmuture of finite collineation groups containing symmetries of generatized quadrangles, Inventiones Math. 40 (1977), 245-265.

## H. A. WILBRINK

## NEARAFFINE PLANES


#### Abstract

In this paper we develop a theory for nearaffine planes analogous to the theory of ordinary affine transtation planes. In a subsequent paper we shall use this theory to give a characterization of a certain class of Minkowski planes.


## 1. Introduction

Nearaffine spaces were introduced by J. André as a generalization of affine spaces (see e.g., [1], [2], [3]). We shall restrict our attention to nearaffine spaces of dimension 2, the nearaffine planes. Our set of axioms, defining nearaffine planes is weaker than the one used by André If, however, the socalled Veblen-condition is assumed to hold (see Section 3), our definition coincides with the one given by Andre in [2]. Our main goal will be to generalize the theory of translation planes to the case of nearaffine planes. In a second paper, we shall show the relationship between certain nearaffine planes and Minkowski planes.

In Section 2 we give the definition of a near affine plane and some basic results. Section 3 is devoted to the so-called Veblen-axiom. In Section 4 we consider automorphisms of nearaffine planes, in particular translations and dilatations. In Section 5 we show that translations exist whenever a certain Desarguers configuration holds. In Section 6 we give an algebraic representation for nearaffine translation planes. Section 7 contains some information on the relationship with Latin squares. Finally, in Section 8, we give a construction of a class of nearaffine planes. More detailed information, especially on the construction of nearaffine planes, can be found in [12].

## 2. Definition and basic results

Let $X$ be a nonempty set of elements called points, $L$ a set of subsets of $X$ called lines. Let $\square$ be an operation called join mapping the ordered pairs $(x, y), x, y \in X, x \neq y$, onto $L$ (the join from $x$ to $y$ is denoted by $x \sqcup y$ ), and $\|$ an equivalence relation called parallelism on $L$ ( $l$ parallel to $m$ is denoted by $l \| m$ ).

We say that $(X, L, L, \|)$ is a nearaffine plane if the following three groups of axioms are satisfied.

Axioms on Lines:

$$
\begin{align*}
& x, y \in x \sqcup y \text { for all } x, y \in X, x \neq y  \tag{L1}\\
& z \in x \sqcup y \backslash\{x\} \Leftrightarrow x \bigsqcup y=x \sqcup z \text { for all } x, y, z \in X, x \neq y \tag{L2}
\end{align*}
$$

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$$
\begin{align*}
& x \bigsqcup y=y \sqcup x=x \bigsqcup z \Rightarrow x \bigsqcup z=z \bigsqcup x \text { for all } x, y, z \in X,  \tag{L3}\\
& y \neq x \neq z:
\end{align*}
$$

The point $x$ is called a basepoint of the line $x \bigsqcup y$. It is not difficult to show the following proposition (see [2]).
PROPOSITION 2.1. The following are equivalent.
(i) $x \sqcup y$ has a basepoint $\neq x$,
(ii) each point of $x \sqcup y$ is a base point of $x \sqcup y$,
(iii) $x \bigsqcup y=y \bigsqcup x$.

Therefore we may define : a line $x \bigsqcup y$ is called straight iff $x \bigsqcup y=y \bigsqcup\rfloor$. The set of all straight lines is denoted by $G$. The lines in $L \backslash G$ are called proper lines.

Axioms of parallelism:
(P1) for all $l \in L, x \in X$ there exists exactly one line with base point $x$ parallel to $l$.
We denote this line by $(x \| l)$.
$x \amalg y \| y \sqcup x \quad$ for all $x, y \in X, x \neq y$.
(P3)
$(g \| l) \Rightarrow l \in G \quad$ for all $g \in G, l \in L$.
Axioms on richness:
(R1) There exists at least two non-parallel straight lines.
(R2) Every line I meets every straight line $g$ with $g \| I$ in exactly one point.
We state some basic results which follow immediately from our axioms (see e.g. [2], [11]).
PROPOSITION 2.2. Two distinct lines with the same base point have no other point in common.
PROPOSITION 2.3. Two distinct straight lines intersect in one point unless they are parallel in which case they are disjoint.
THEOREM 2.4. A nearaffine plane with commutative join is an affine plane.
We shall only consider finite nearaffine planes, i.e., nearaffine planes with a finite number of points. The following result is easy to prove (see, e.g., [2], [11]).
PROPOSITION 2.5. All lines of a nearaffine plane have the same number of points.

The number of points on a line, which equals the number of parallel straight lines in one equivalence class, is denoted by $n$ and called the order of the nearaffine plane.

## NEARAFFINE PLANES

PROPOSITION 2.6. $|X|=n^{2}$.
PROPOSITION 2.7. There are exactly $n+1$ lines with a giten base point.
We denote by $s+1$ the number of equivalence classes containing straight lines. $\mathrm{By}(\mathrm{RI})$ we have $s \geqslant 1$.
PROPOSITION 2.8. Every point is on $s+1$ straight lines, $|G|=n(s+1)$, $|L G|=n^{2}(n-s)$.

## 3. The veblen-condition

Many interesting examples of nearaffine planes le.g., the nearaffine planes associated with Minkowski planes) satisfy the following version of the Veblen-condition (named ( $V^{\prime}$ ) in [2]).
( $V^{\prime}$ ) Let $q$ be a straight line, $P, Q, R$ distinct points on $g, I \neq g$ a line with base point $P$ and $S \in \Lambda\{P$. Then $(R \| Q \sqcup S) \cap \mid \neq \varnothing$ (see Figure 1).


Fig 1
Before we prove the main result on nearaffine planes which satisfy ( $V^{\prime}$ ), we prove a proposition valid in any nearaffine plane. Notice that until now we have not used axiom ( P 2 ) and that the proof of this proposition only requires the following weakened version of ( P 2 ) (this will be important in our paper on Minkowski planes).
( $\mathbf{P}^{\prime}$ ) Let $g$ and $h$ be two distinct parallel straight lines, $x, x^{\prime} \in g$ and fo $y^{\prime} \in h$. Then $x \square y\left\|x^{\prime} \bigsqcup y^{\prime} \Leftrightarrow y^{\prime} \square x\right\| y^{\prime} \square x^{\prime}$.

PROPOSITION 3.1. Two parallel lines which have their base point on one straight line are disjoint or identical.

Proof. Let $/$ and $l^{\prime}$ be two parallel lines with base points $x$ and $x^{\prime}$ respectively on the straight line $g$. If $y \in I \cap l^{\prime}, y \neq x, x^{\prime}$ then $x \sqcup y=I \| I^{\prime}=x^{\prime} \bigsqcup y$, hence $y \sqcup x \| y \sqcup x^{\prime}$ by $\left(\mathbf{P}^{\prime}\right)$ and so $y \sqcup x=y \amalg x^{\prime}$ by (PI). Therefore $x=x^{\prime}$ by ( R 2 ) and so $l=l^{\prime}$ by ( P 1 ).

THEOREM 3.2. (André [2]). Let $1=(X, L, L, \mid)$ be a nearaffine plane satisfing ( $V^{\prime}$ ) and $g$ a straight line of 1 . Then the point set $X$ and the line set $L_{g}:=\{l \in L| |$ has base point on $g\} \cup\{\in G \mid h \| g\}$ constitute an affine plane $\mathrm{j}_{g}^{\prime}=\left(X, L_{i}\right)$.
Proof. Let $l, m \in L_{g}, l \neq m$. If $l \mid m$ then $|I \cap m|=0$ by 2.3 and 3.1 . If $I \mid f$
then $|l \cap m|=1$. This follows from (R2) if $l \| g$ or $m \| g$. Suppose, therefore, that $l$ and $m$ have base points on $g$. The $n$ line in $L_{g}$ parallel to $m$ partition $X$ by 3.1. Hence, at least one of these lines contains a point of $l$. Therefore, by $\left(V^{\prime}\right)$ and 2.5 , each of these lines, so in particular $m$, contains exactly one point of $l$. Since $\left|L_{g}\right|=n(n+1)$ and $|l|=n$ for every $l \in L_{g}$ it follows from [5, result 3.2.4c, p. 139] that $\mathscr{N}_{g}$ is an affine plane.
Remark. Notice that two lines of $\mathscr{N}_{g}$ are parallel in $\mathscr{N}_{g}$ (i.e., disjoint) iff they are parallel in $\mathscr{A}$.

## 4. Automorphisms

In this section we generalize such notions as automorphism, dilatation etc. to the case of nearaffine planes. Proofs which do not differ essentially from the corresponding proofs for affine planes (see e.g., [4]) will be omitted.
DEFINITION 4.1. Let $\mathscr{N}=(X, L, L, \|)$ and $\mathscr{N}^{\prime}=\left(X^{\prime}, L^{\prime}, L^{\prime}, \|^{\prime}\right)$ be two nearaffine planes. A bijection $\alpha: X \rightarrow X^{\prime}$ is called an isomorphism of $\mathcal{N}$ and $\mathcal{N}^{\prime}$ if
(i) $(P \sqcup Q)^{\alpha}=P^{\alpha} \sqcup^{\gamma} Q^{\alpha} \quad$ for all $P, Q \in X, P \neq Q$, and
(ii) $l\left\|m \Leftrightarrow l^{\alpha}\right\|^{\prime} m^{\alpha}$ for all $l, m \in L$.

If $\mathscr{N}=\mathscr{N}^{\prime}$, then $\alpha$ is called an automorphism of $\mathscr{N}$. A permutation $\alpha$ of the points of $\mathscr{N}$ is called a dilatation if $P \sqcup Q \| P^{\alpha} \sqcup Q^{\alpha}$ for all $P \neq Q$.

The automorphisms of a nearaffine plane form a group $\mathscr{A}$, the dilatations form a group $\mathscr{D}$.

THEOREM 4.2. $\mathscr{D} \leqslant \mathscr{A}$.
LEMMA 4.3. Suppose $\delta \in \mathscr{D}$ fixes $P \in X$. Then $Q^{\delta} \in P \sqcup Q$ for all $Q \in X$, $X \neq P$.

THEOREM 4.4. Suppose $\delta \in \mathscr{D}$ fixes two distinct points $P$ and $Q$. Then $\delta=1$.
Proof. Take $R \in X$. If $R=P$ or $R=Q$, then $R^{\delta}=R$. if $R \neq P, Q$ we have by 4.3: $R^{\delta} \in P \bigsqcup R$ and $R^{\delta} \in Q \sqcup R$. By (R1) there is at least one straight line $g \neq P \sqcup Q$ through $P$, so for $R \in g$ we have $R^{\delta} \in(P \sqcup R) \cap(Q \sqcup R)=\{R\}$, i.e., $R^{\delta}=R$. For an arbitrary $R \notin g$ we replace $P$ by a point $P^{\prime}$ in such a way that $P^{\prime} \sqcup R$ is straight and $Q$ by some point $Q \in g \backslash\left\{P^{\prime}\right\}$. It follows that $R^{\delta} \in\left(P^{\prime} \mid \sqcup R\right)$ $\cap\left(Q^{\prime} \sqcup R\right)=\{R\}$.
COROLLARY 4.5. Let $\delta_{1}, \delta_{2} \in \mathscr{D}$ and suppose $P^{\delta_{1}}=P^{\delta_{2}}, Q^{\delta_{1}}=Q^{\delta_{2}}$ for distinct points $P$ and $Q$. Then $\delta_{1}=\delta_{2}$.

DEFINITION 4.6. A dilatation $\tau$ is called a translation if $\tau=1$ or if $P \sqcup P^{t} \| Q \sqcup Q^{\imath}$ for all $P, Q \in X$. The parallel class containing $P \sqcup P^{t}$ is called the direction of $\tau \neq 1$. The translation $\tau$ is straight if $P \bigsqcup P^{\tau}$ is straight. We denote by $\mathscr{F}$ the set of all translations.

A translation $\tau \neq 1$ has no fixd point. Suppose $P^{r}=P$; then for any point $Q \neq P$ we have $Q^{t} \neq Q$ by 4.4 and $Q^{t} \in P \sqcup Q$ by 4.3. Hence, if $P \sqcup Q$ is straight, $Q \sqcup Q^{t}=P \sqcup Q$.

This is a contradiction since there are at least two nonparallel straight lines through $P$.
LEMMA 4.7. If $\alpha \in \mathscr{A}$ and $\tau \in \mathscr{T}$, then $\alpha \tau \alpha^{-1} \in \mathscr{T}$. If in addition $\alpha \in \mathscr{D}$ and $\tau \neq 1$, then $\tau$ and $\alpha \tau \alpha^{-1}$ have the same direction.
THEOREM 4.8. Let $C$ be a parallel class consisting of straight lines and $\mathscr{T}(C):=\{\tau \in \mathscr{F} \mid \tau$ has direction $C\} \cup\{1\}$. Then $\mathscr{F}(C) \& D$.
LEMMA 4.9. Let $C$ and $D$ be two distinct parallel classes consisting of straight lines. Then $\sigma \tau=\tau \sigma$ for all $\sigma \in \mathscr{F}(C), \tau \in \mathscr{T}(D)$.
LEMMA 4.10. Let $C$ and $D$ be two parallel classes containing straight lines, $\sigma \in \mathscr{F}(C)$ and $\tau \in \mathscr{T}(D)$. If $\sigma \tau \neq 1$, then $\sigma \tau$ has no fixed points.
Proof. If $C=D$ or if $\sigma$ or $\tau=1$, this is a consequence of 4.8. If $C \neq D$ and $\sigma, \tau \neq 1$, then $P^{\sigma \tau}=P$ for some $P \in X$ implies $P \bigsqcup P^{\sigma} \in C, P \bigsqcup P^{\tau-1} \in D$, $P^{\sigma}=P^{r^{-1}}$, a contradiction.

For nearaffine planes the product of two translations need not be a translation. For straight translations the following theorem holds.
THEOREM 4.11. Let $C, D$ and $E$ be three distinct parallel classes consisting of straight lines. Suppose $\rho \in \mathscr{T}(C), \sigma \in \mathscr{T}(D), \tau \in \mathscr{T}(E)$ and $P \in X$ satisfy $P^{\rho \sigma}=P^{\top}$. Then $\rho \sigma=\tau$.

Proof. If $\tau=1$, then $P^{\rho \sigma}=P$, hence $\rho \sigma=1$ by 4.10 . If $\tau \neq 1$, then $P^{\mathrm{r}} \neq P$. From 4.9 it follows that $\left(P^{t}\right)^{\tau}=\left(P^{\rho \sigma}\right)^{\tau}=\left(P^{\tau}\right)^{\rho \sigma}$. Hence, $\tau=\rho \sigma$ by 4.5 .
THEOREM 4.12. Let $C$ and $D$ be two distinct parallel classes consisting of straight lines with $|\mathscr{T}(C)|=|\mathscr{F}(D)|=n$. Then

$$
\mathscr{T} \subseteq\langle\mathscr{T}(C), \mathscr{T}(D)\rangle=\mathscr{T}(C) \mathscr{T}(D)
$$

If in addition $\mathscr{T}(C)$ and $\mathscr{T}(D)$ are Abelian, then $\mathscr{T}=\mathscr{T}(C) \mathscr{T}(D)$.
Proof. By $4.9,\langle\mathscr{F}(C), \mathscr{T}(D)\rangle=\mathscr{F}(C) \mathscr{T}(D)$ and $|\mathscr{T}(C) \mathscr{F}(D)|=n^{2}$. By 4.10, $\mathscr{F}(C) \mathscr{T}(D)$ is the Frobenius kernel of $\mathscr{D}$, hence it contains all fixed-points free dilatations. Therefore $\mathscr{T} \subseteq \mathscr{F}(C) \mathscr{T}(D)$. Suppose $\mathscr{T}(C)$ and $\mathscr{T}(D)$ are Abelian. Take $\rho \in \mathscr{F}(C), \sigma \in \mathscr{T}(D)$ and $P, Q \in X$. There exist $\rho_{1} \in \mathscr{T}(C), \sigma_{1} \in \mathscr{F}(D)$ such that $P^{\rho \mid \sigma_{1}}=Q$. Hence,

$$
P \bigsqcup P^{\rho \sigma} \|\left(P \sqcup P^{\rho \sigma}\right)^{\rho_{1} \sigma_{1}}=P^{\rho_{1} \sigma_{1}} \bigsqcup P^{\rho_{1} \sigma_{1} \rho \sigma}=Q \sqcup Q^{\rho \sigma},
$$

i.e., $\rho \sigma \in \mathscr{T}$.

A nearaffine plane having two distinct parallel classes $C$ and $D$ consisting of straight lines such that $|\mathscr{F}(C)|=|\mathscr{T}(D)|=n$ is called a nearaffine translation plane. Notice that this definition is consistent with the definition of translation plane.

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THEOREM 4.13. Let $C, D$ and $E$ be three distinct parallel classes consisting of straight lines. $I f|\mathscr{F}(C)|=|\mathscr{T}(D)|=n$, then
(a) $\mathscr{T}(E)$ is Abelian,
(b) $\mathscr{F}(C) \simeq \mathscr{T}(D)$.

Proof. (a) Let $\tau_{1}, \tau_{2} \in \mathscr{F}(E)$. By 4.12 there exist $\rho_{1} \in \mathscr{T}(C), \sigma_{1} \in \mathscr{T}(D)$ such that $\tau_{1}=\rho_{1} \sigma_{1}$. By 4.9,

$$
\tau_{1} \tau_{2}=\rho_{1} \sigma_{1} \tau_{2}=\tau_{2} \rho_{1} \sigma_{1}=\tau_{2} \tau_{1} .
$$

(b) Define the automorphism $\phi: \mathscr{T}(C) \rightarrow \mathscr{T}(D)$ as follows: Fix a line $g \in E$. For each $\rho \in \mathscr{T}(C)$ let $\phi(\rho) \in \mathscr{T}(D)$ be determined by $g^{\rho \phi(\rho)}=g$.
COROLLARY 4.14. If in addition to the hypothesis of $4.13,|T(E)|=n$, then $\mathscr{T}(C) \simeq \mathscr{T}(D) \simeq \mathscr{F}(E)$ and these groups are Abelian.

So far we have not used ( P 2 ) in this section. Using ( P 2 ) it is possible to prove the following theorem.
THEOREM 4.15. The order $n$ of a nearaffine translation plane is odd or a power of 2 .

Proof. Suppose $n$ is even and let $C$ and $D$ be two distinct parallel classes consisting of straight lines such that $|\mathscr{T}(C)|=|\mathscr{T}(D)|=n$. There exists $\rho \in \mathscr{F}(C)$ such that $\rho^{2}=1, \rho \neq 1$. Take $\sigma \in \mathscr{F}(D)$ and $P \in X$. Then,

$$
P \bigsqcup P^{\rho \sigma}\left\|P^{\rho \sigma-1} \bigsqcup\left(P^{\rho \sigma}\right)^{\rho \sigma-1}=P^{\rho \sigma-1} \bigsqcup P\right\| P \bigsqcup P^{\rho \sigma-1} .
$$

Therefore $P^{\rho \sigma}, P^{\rho \sigma-1} \in P \sqcup P^{\rho \sigma} \notin D$. Since $P^{\rho \sigma-1}$ and $P^{\rho \sigma}=\left(P^{\rho \sigma-1}\right)^{\sigma 2}$ are on the same straight line of $D$ it follows that $P^{\rho \sigma-1}=P^{\rho \sigma}$, i.e., $\sigma^{2}=1$. Hence. $\mathscr{F}(D)$ is an (elementary Abelian) 2-group.

## 5. A desargues configuration

Let. $\mathcal{N}=(X, L, \amalg, \|)$ be a nearaffine plane and $C$ a parallel class consisting of straight lines. Consider the following condition (cf. [2], [3]).
(D1) Little Desargues configuration. If $P, P^{\prime}, Q, Q^{\prime}, R, R^{\prime} \in X$ are distinct points such that $P \sqcup P^{\prime}, Q \sqcup Q, R \bigsqcup R^{\prime}$ are distinct lines of $C$, then $P \sqcup Q \| P^{\prime} \sqcup Q^{\prime}$ and $P \sqcup R \| P^{\prime} \sqcup R^{\prime}$ imply $Q \sqcup R \| Q^{\prime} \sqcup R^{\prime}$ (see Figure 2).

Analogous to the situation for affine planes, the validity of (D1) is seen to be equivalent to the existence of all possible translations with direction $C$.
THEOREM 5.1. C satisfies (D1) $\Leftrightarrow|\mathscr{T}(C)|=n$.
The following theorem will be useful in our paper on Minkowski planes. Again notice that we only make use of ( $\mathbf{P} \mathbf{2}^{\prime}$ ).

THEOREM 5.2. Let $\mathscr{N}=(X, L, \sqcup, \|)$ be a nearaffine plane in which the Veblen-condition holds, and let $C$ be a parallel class of straight lines. Then

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Fig. 2.
(using the notation of 3.2 ), $C$ satisfles (D1) in $\mathcal{N} \Leftrightarrow C$ satisfies (D1) in $\mathcal{N}_{g}$ for all $g \in C$.

Proof. $\Rightarrow$ : Every translation of $A$ with direction $C$ is easily seen to induce a translation of $A_{g}$, with direction $C$ for every $g \in G$.
$\Leftarrow:$ Let $P, P^{\prime}, Q, Q^{\prime}, R, R^{\prime}$ be distinct points such that $P \bigsqcup P^{\prime}, Q \sqcup Q$, and $R \sqcup R^{\prime}$ are distinct straight lines of $C$ and such that $P \sqcup Q \| P^{\prime} L Q^{\prime}$. $P \sqcup R \| P^{\prime} \sqcup R^{\prime}$. Let $S$ (resp. $S^{\prime}$ ) be the base point of the line in $\mathcal{N}_{P \sqcup P^{\prime}}$ passing through $Q$ and $R$ (resp. $Q^{\prime}$ and $R^{\prime}$ ), (see Figure 3). Application of (D1) in $\mathscr{N}_{p \sqcup p}$. yields $S \bigsqcup Q \| S^{\prime} \bigsqcup^{\prime} Q^{\prime}$.


Fig. 3.
Let $D$ be a parallel class of straight lines different from $C$, and let $T$ (resp. $T^{\prime}$ ) be the point of intersection of $P \bigsqcup P^{\prime}$ and the straight line of $D$ passing through $R$ (resp. $R^{\prime}$ ). Application of (D1) in $\mathcal{A}_{P \sqcup P^{\prime}}$ to the triangles $T Q R$ and $T^{\prime} Q^{\prime} R^{\prime}$ yields $T \sqcup Q \| T^{\prime} \sqcup Q^{\prime}$, hence $Q \sqcup T \| Q^{\prime} \sqcup T^{\prime}$. Finally apply (D1) in $\mathcal{V}_{Q \cup Q^{\prime}}^{\prime}$ to the triangle $T Q R$ and $T^{\prime} Q^{\prime} R^{\prime}$ to obtain $Q \sqcup R \| Q^{\prime} \sqcup R^{\prime}$.

## 6. Algebraic representation

In this section an algebraic representation is given of the nearaffine translation planes. The tedious but straightforward proofs are omitted. For details see [12].

Let $G$ and $G^{\prime}$ be two groups of order $n$ written additively. We do not assume that $G$ or $G^{\prime}$ is Abelian or that $G \simeq G^{\prime}$ (although the same symbol + is used

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for addition in both groups). Let $\mathscr{F}$ be a set of $(n-1)$ mappings $f_{i}: G \rightarrow G^{\prime}$, $i=1, \ldots, n-1$, such that the following conditions are satisfied.
(i) $f_{i}$ is a bijection for all $i=1, \ldots, n-1$.
(ii) $f_{i}(0)=0$ for all $i=1, \ldots, n-1$.
(iii) $f_{i}(\alpha)=-f_{i}(-\alpha)$ for all $i=1, \ldots, n-1, \alpha \in G$.
(iv) $f_{i}(\alpha) \neq f_{j}(\alpha)$ for $1 \leqslant i<j \leqslant n-1, \alpha \in G \backslash\{0\}$.
(v) For all $i=1, \ldots, n-1$ either,

$$
\forall_{\alpha \in G \backslash(0)} \exists_{\beta \in G}\left[f_{i}(\alpha+\beta) \neq f_{i}(\alpha)+f_{i}(\beta)\right]
$$

or

$$
\forall_{\alpha, \beta \in \mathrm{G}}\left[f_{i}(\alpha+\beta)=f_{i}(\alpha)+f_{i}(\beta)\right]
$$

and $f_{i}-f_{j}$ is a bijection for $j=1, \ldots, n-1, j \neq i$.
Given such a set of mappings $\mathscr{F}$ it is possible to construct a nearaffine translation plane in the following way. Put $X:=G \times G^{\prime}$. For $x, y \in X$, $x=\left(\xi, \xi^{\prime}\right), g=\left(\eta, \eta^{\prime}\right), x \neq y$, define:

$$
x \sqcup y:= \begin{cases}\left\{\left(\xi, \alpha^{\prime}\right) \mid \alpha^{\prime} \in G^{\prime}\right\} & \text { if } \xi=\eta, \\ \left.\left(\alpha, \xi^{\prime}\right) \mid \alpha \in G\right\} & \text { if } \xi^{\prime}=\eta^{\prime}, \\ \left\{\left(\xi+\alpha, \xi^{\prime}+f_{i}(\alpha) \mid \alpha \in G\right\}\right. & \text { if } \xi \neq \eta, \xi^{\prime} \neq \eta^{\prime} \text { and } \\ & f_{i}(-\xi+\eta)=-\xi^{\prime}+\eta^{\prime} .\end{cases}
$$

The line set $L$ is just the set of all $x \sqcup y, x \neq y$. For any line $l=x \sqcup y$ we let $\mathrm{d}(l) \in\{0,1, \ldots, n-1, \infty\}$ be determined by

$$
\mathrm{d}(l):= \begin{cases}\infty & \text { if } \xi=\eta, \\ 0 & \text { if } \xi^{\prime}=\eta^{\prime}, \\ \mathrm{i} & \text { if } \xi^{\prime} \neq \eta^{\prime}, \xi^{\prime} \neq \eta^{\prime} \text { and } f_{i}(-\xi+\eta)=-\xi^{\prime}+\eta^{\prime} .\end{cases}
$$

Notice that $\mathrm{d}(l)$ only depends on $l$ and not on the special choice of $x$ and $y$. Define parallelism by

$$
l \| m: \Leftrightarrow \mathrm{d}(l)=\mathrm{d}(m)
$$

then $\mathscr{N}=(X, L, \amalg, \|)$ is a nearaffine translation plane. Conversely, every nearaffine translation plane can be described in this way. The parallel classes $C_{0}:=\{l \in L \mid \mathrm{d}(l)=0\}, C_{\infty}:=\{l \in L \mid \mathrm{d}(l)=\infty\}$ consist of straight lines. For each $\alpha \in G$, the mapping $\left(\xi, \xi^{\prime}\right) \rightarrow\left(\alpha+\xi, \xi^{\prime}\right)$ is a translation with direction $C_{0}$. For each $\alpha^{\prime} \in G^{\prime}$, the mapping $\left(\xi, \xi^{\prime}\right) \rightarrow\left(\xi, \alpha^{\prime}+\xi^{\prime}\right)$ is a translation with direction $C_{\infty}$. For $i=1, \ldots, n-1, C_{i}:=\{l \in L \mid \mathrm{d}(l)=i\}$ consists of straight lines iff $i$ satisfies the second alternative of (v).

The Veblen-condition $\left(V^{\prime}\right)$ is satisfied if for $1 \leqslant i<j \leqslant n-1$,
(a) $f_{i}-f_{i}: G \rightarrow G^{\prime}$ is a bijection,
(b) $\hat{f}_{i}-f_{j}: G^{\prime} \rightarrow G$ is a bijection,
(c) for all $k \in\{1, \ldots, n-1\}$ which satisfy the second alternative of $(v)$ and for all $\gamma \in G$ there is a unique solution $\alpha$ of $f_{k}(\gamma)=f_{j}(\gamma+\alpha)-f_{i}(\alpha)$.

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## 7. Nearaffine planes and latin squares

It is well known that the existence of an affine plane of order $n$ is equivalent to the existence of $n-1$ mutually orthogonal Latin squares (MOLS) of order $n$ (see [5]). For nearaffine plane the following result holds.

THEOREM 7.1. If $\mathscr{N}$ is a nearaffine plane of order $n$ with $s+1$ parallel classes containing straight lines ( $s<n$ ), then there exist $s$ MOLS of order $n$.

Proof. The $n(s+1)$ lines in the $s+1$ parallel classes consisting of straight lines together with $n$ lines from a parallel class consisting of proper lines, all having their base points on a fixed straight line, constitute an $(s+2)-$ net of order $n$. This is equivalent to the existence of $s$ MOLS of order $n$ (see, e.g., [5]).

Let $N$ be an integer, $N \geqslant 2$, and suppose $N$ has prime decomposition $N=$ $=p_{1}^{\alpha_{1}} p_{2}^{z_{2}} \ldots p_{k}^{\alpha_{k}}$. Define

$$
s(N):=\min _{1 \leqslant i \leqslant k} p_{i}^{\alpha_{i}}-1 .
$$

It is well known (see, e.g., [5]) that there exist at least $s(N)$ MOLS of order $N$ (the so-called MacNeish bound). The following theorem shows therefore that, as far as nearaffine translation planes are concerned, we cannot hope for interesting applications of 7.1

THEOREM 7.2. Let $\mathcal{N}$ be a nearaffine translation plane of order $n$ with $s+1$ parallel classes containing straight lines. Then $s \leqslant s(n)$.

Proof. Notice that $s-1$ of the $f_{i}^{\prime}$ s associated with $\mathcal{N}$, say $f_{1}, f_{2}, \ldots, f_{s-1}$, satisfy the second alternative of (v) of Section 6. Put $\phi_{i}:=\bar{f}_{1} \circ f_{i}, i=1,2, \ldots, s$. Then $\phi_{i}-\phi_{j}: G \rightarrow G$ is a permutation of the elements of $G, 1 \leqslant i<j \leqslant s$. Hence, the Latin squares $A^{(i)}=\left[a_{x, y}^{(i)}\right]$ defined by

$$
a_{x, y}^{(i)}:=\phi_{i}(x)+y, \quad i=1, \ldots, s, \quad x, y \in G,
$$

are mutually orthogonal. Since $\phi_{1}, \phi_{2}, \ldots, \phi_{s-1}$ are automorphisms of $G$ it follows by a theorem of H. B. Mann (see [6] or [9]) that $s-1 \leqslant s(n)$. Suppose $s-1=s(n)=p^{\alpha}-1, p$ a prime, $\alpha \in \mathbb{N}$. It follows from the proof of Mann's theorem that the elements $\neq 0$ of a Sylow $p$-subgroup $P$ of $G$ are all in different conjugacy classes. Thus $-y+x+y \in P \Rightarrow-y+x+y=x$ for all $x \in P, y \in G$. In particular, if $y \in N_{G}(P)$ then $x+y=y+x$ for all $x \in P$, i.e., $P \leqslant Z\left(N_{G}(P)\right)$. By a theorem of Burnside (see [7] or [8]), $G$ contains a normal $p$-complement $N$. Since $|G \backslash N|$ and $|N|$ are coprime, $N$ is a characteristic subgroup of $G$. Thus, the rows and columns of $A^{(1)}, \ldots, A^{(s-1)}$ which correspond to the elements of $N$, yield mutually orthogonal Latin subsquares of order $n / p^{\alpha}$. By a theorem of Parker (see [10]) such a set of $s-1$ MOLS cannot be extended to a set of $s$ MOLS, a contradiction. Hence, $s-1<s(n)$, i.e., $s \leqslant s(n)$.

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## 8. Construction of nearaffine planes

Using the representation of nearaffine translation planes of Section 6, we treat a special case of the more general construction described in [12]. The nearaffine planes thus obtained turn out to be associated with certain Minkowski planes. Let $p$ be a prime, $h$ a positive integer and $n=p^{h}$. For the groups $G$ and $G^{\prime}$ of Section 6 we take the additive group of $G F(n)$. Fix an automorphism $\phi$ of $G F(n)$, and for each $a \in G F(n)^{*}$ define $f_{a}: G F(n) \rightarrow G F(n)$ by $f_{a}(0):=0$ and

$$
\begin{aligned}
& f_{a}(x):=a x^{-1}, \quad x \in G F(n)^{*}, \text { if } a \text { is a square, } \\
& f_{a}(x):=a\left(x^{-1}\right)^{\phi}, \quad x \in G F(n)^{*}, \text { if } a \text { is a nonsquare. }
\end{aligned}
$$

The set $\mathscr{F}:=\left\{f_{a} \mid a \in G F(n)^{*}\right\}$ is easily seen to satisfy the properties (i),...,(v) of Section 6. The corresponding nearaffine plane is of order $n$, and $s=1$. It is also not hard to show that the Veblen-condition holds in these nearaffine planes.

## REFERENCES

1. André, J.- Eine Kennzeichung der Dilatationsgruppen desarguesscher affiner Räume als Permutationsgruppen', Arch. Math. 25 (1974), 411-418.
2. André, J.: 'Some New Results on Incidence Structures', Atti dei Convegni Lincei 17, Colloquio Internazionale sulle teorie combinatori II (1976), pp. 201-222.
3. André. J. : 'On Finite Non-commutative Affine Spaces’. in M. Hall and J. H. von Lint (eds.) Combinatorics Part I, Mathematical Centre Tracts 55 (1974), 60-107.
4. Artin, E. : Geometric Algebra, Interscience, New York, London, 1957.
5. Dembowski, P., Finite Geometries, Springer-Verlag, Berlin, Heidelberg, New York, 1968.
6. Denes, J. and Keedwell, A. D.: Latin Squares and their Applications, The English Universities Press Ltd., London, 1974.
7. Hall, M.: The Theory of Groups, MacMillan, New York, 1959.
8. Huppert, B.: Endliche Gruppen I, Springer-Verlag, Berlin, Heidelberg, New York, 1967.
9. Mann, H. B.: 'The Construction of Orthogonal Latin-Squares', Ann. Math. Stat. (1943), 401-414.
10. Parker, E. T.: Nonextendibility Conditions on Mutually Orthogonal Latin-Squares", Proc. Amer. Math. Soc. 13 (1963), 219-221.
11. Van der Schoot, J. and Wilbrink, H.: 'Nearaffine Planes I', Indag. Math. 37 (1975), 137-143.
12. Wilbrink, H.: 'Nearaffine Planes and Minkowski Planes', Master's thesis, Techn. Univ. Eindhoven, 1978.

## FINITE MINKOWSKI PLANES

Abstract. In this paper we give second characterizations of a certain class of finite Minkowski planes.

## 1. Introduction

It is well known, see e.g. [5], that with each point of a Minkowski plane there is associated an affine plane, its so-called derived plane. It is the purpose of this paper to show that, under certain additional hypotheses, with each point of a Minkowski plane there is also associated a nearaffine plane, its residual plane. In addition we show that the 'known' Minkowski plane are characterized by the fact that these nearaffine planes are nearaffine translation planes (see [9]). Using this result a configurational condition is obtained in a completely natural way which characterizes the known Minkowski planes.

## 2. BASIC CONCEPTS

Let $M$ be a set of points and $\mathscr{L}^{+}, \mathscr{L}^{-}, \mathscr{C}$ three collections of subsets of $M$. The elements of $\mathscr{L}:=\mathscr{L}^{+} \cup \mathscr{L}^{-}$are called lines or generators, the elements of $\mathscr{C}$ are called circles. We say that $\mathscr{M}=\left(M, \mathscr{L}^{+}, \mathscr{L}^{-}, \mathscr{C}\right)$ is a Minkowski plane if the following axioms are satisfied (cf. [5]):
(M1): $\quad \mathscr{L}^{+}$and $\mathscr{L}^{-}$are partitions of $M$.
(M2): $\quad\left|l^{+} \cap l^{-}\right|=1 \quad$ for all $l^{+} \in \mathscr{L}^{+}, l^{-} \in \mathscr{L}^{-}$.
(M3): Given any three points no two on a line, there is a unique circle passing through these three points.
(M4): $\quad|\ln c|=1 \quad$ for all $l \in \mathscr{L}, c \in \mathscr{C}$.
(M5): There exist three points no two of which are on one line.
(M6): $\quad$ Given a circle $c$, a point $P \in c$ and $a$ point $Q \notin c, P$ and $Q$ not on one line, there is a unique circle $d$ such that $P, Q \in d$ and $c \cap d=$ $=\{P\}$.
Two points $P$ and $Q$ are called plus-parallel (notation $P \|_{+} Q$ ) if $P$ and $Q$ are on a line of $\mathscr{L}^{+}$, minus-parallel $\left(P \|_{-} Q\right)$ if $P$ and $Q$ are on a line of $\mathscr{L}^{-}$. Parallel $(P \| Q)$ means either $P \|_{+} Q$ or $P \|_{-} Q$. For $P \in M, \epsilon=+,-$ we denote by $[P]_{\epsilon}$ the unique line in $\mathscr{L}^{\epsilon}$ incident with $P$. If $P, Q$ and $R$ are (distinct) nonparallel points, then we denote by $(P, Q, R)$ the unique circle containing $P, Q$ and $R$. Two circles $c$ and $d$ touch in a point $P$ if $c \cap d=\{P\}$.

Fix a point $Z$ and put

$$
\begin{aligned}
& M_{Z}:=M \backslash\left([Z]_{+} \cup[Z]_{-}\right), \\
& L_{Z}:=\left\{c^{*} \mid c \in \mathscr{C}, Z \in c\right\} \cup\left\{l^{*} \mid l \in \mathscr{L} \backslash\left\{[Z]_{+},[Z]_{-}\right\}\right\}
\end{aligned}
$$

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where the * indicates that we have removed the point that the circle or line has in common with $[Z]_{+} \cup[Z]_{-}$. Then $\mathscr{M}_{z}:=\left(M_{Z}, L_{z}\right)$ is an affine plane with pointset $M_{z}$ and lineset $L_{z}$ (see, e.g., [5]). We call $\mathscr{M}_{z}$ the derived plane with respect to the point $Z$. We shall only consider finite Minkowski planes, i.e., Minkowski planes with a finite number of points. For finite Minkowski planes (M6) is a consequence of the other axioms (see [5]). It is easily seen that $\left|\mathscr{L}^{+}\right|=\left|\mathscr{L}^{-}\right|=|l|=|c|=: n+1$ for all $l \in \mathscr{L}, c \in \mathscr{C}$. The integer $n$ is called the order of the Minkowski plane. Notice that $n$ is also the order of the derived planes $\mathscr{M}_{z}$.

Following Benz [1] we sketch the close relationship between (finite) Minkowski planes and sharply 3 -transitive sets of permutations. Let $\Omega$ be a finite set, $|\Omega|=n+1 \geqslant 3$, and $G$ a subset of $S^{\Omega}$, the symmetric group on $\Omega$, acting sharply triply transitively on $\Omega$.

Define

$$
\begin{aligned}
M & :=\Omega \times \Omega \\
\mathscr{L}^{+} & :=\{\{(\alpha, \beta) \mid \alpha \in \Omega\} \mid \beta \in \Omega\} \\
\mathscr{L}^{-} & :=\{\{(\alpha, \beta) \mid \beta \in \Omega\} \mid \alpha \in \Omega\}, \\
\mathscr{C} & :=\left\{\left\{\left(a, \alpha^{g}\right) \mid \alpha \in \Omega\right\} \mid g \in G\right\} .
\end{aligned}
$$

Then $\mathscr{M}:=(\Omega, G):=\left(M, \mathscr{L}^{+}, \mathscr{L}^{-}, \mathscr{C}\right)$ is a Minkowski plane of order $n$. Conversely, every Minkowski plane can be obtained in this way.

Two Minkowski planes $\mathscr{M}=(\Omega, G)=\left(M, \mathscr{L}^{+}, \mathscr{L}^{-}, \mathscr{C}\right)$ and $\mathscr{M}^{\prime}=\left(\Omega^{\prime}, G^{\prime}\right)=$ $=\left(M^{\prime}, \mathscr{L}^{+\prime}, \mathscr{L}^{-\prime}, \mathscr{C}^{\prime}\right)$ are said to be isomorphic if there is a bijection $s: M \rightarrow M^{\prime}$ such that

$$
\mathscr{L}^{s}=\mathscr{L}^{\prime} \quad \text { and } \quad \mathscr{C}^{s}=\mathscr{C}^{\prime} .
$$

Since $s$ maps the disjoint lines of $\mathscr{L}^{+}$onto disjoint lines there are only two possibilities, either $\left(\mathscr{L}^{q}\right)^{s}=\mathscr{L}^{z}$ or $\left(\mathscr{L}^{e}\right)^{s}=\mathscr{L}^{-\varepsilon}, \varepsilon=+,-$. In the first case $s$ is called a positive isomorphism in the second case a negative isomorphism. If $s$ is a positive isomorphism then there exist bijections $a, b: \Omega \rightarrow \Omega^{\prime}$ such that $(\alpha, \beta)^{s}=\left(\alpha^{a}, \beta^{b}\right)$ for all $\alpha, \beta \in \Omega$, and $G^{\prime}=a^{-1} G b$. If $s$ is a negative isomorphism then there exist bijections $a, b: \Omega \rightarrow \boldsymbol{\Omega}^{\prime}$ such that $(\alpha, \beta)^{s}=\left(\beta^{b}, \alpha^{\alpha}\right)$, and $G^{\prime}=$ $=b^{-1} G^{-1} a$. It follows that we may assume w.l.o.g that id $\in G$.

A (positive, negative) automorphism of a Minkowski plane $\mathscr{M}$ is a (positive, negative) isomorphism of $\mathscr{M}$ onto itself. The automorphism group Aut $(\Omega, G) \leqslant S^{\Omega \times \Omega}$ of the Minkowski plane $(\Omega, G)$ is given by

$$
\operatorname{Aut}(\Omega, G)=\left\{(a, b) \mid a^{-1} G b=G\right\} \cup\left\{(a, b) \mid a^{-1} G b=G^{-1}\right\} \tau
$$

where $\tau$ is the permutation which sends $(\alpha, \beta)$ to $(\beta, \alpha)$.

## 3. The residual plane

Let $\mathscr{M}=\left(M, \mathscr{L}^{+}, \mathscr{L}^{-}, \mathscr{C}\right)$ be a Minkowski plane. Fix a point $Z \in M$ and define $\left.M_{z}=M \backslash[Z]_{+} \cup[Z]_{-}\right)$. We have already remarked that the lines $\neq$
$[Z]_{+},[Z]_{-}$together with the circles which are incident with $Z$ are the lines of an affine plane with pointset $M_{Z}$. We shall show that the lines $\neq[Z]_{+}$, $[Z]_{-}$together with the circles not incident with $Z$ are the lines of a nearaffine plane with the same pointset if suitable conditions are assumed to hold in $\mathscr{M}$.

For each point $P \in M_{Z}$ we let the points $P^{+}$and $P^{-}$be defined by $P^{+}:=$ $[Z]_{+} \cap[P]_{-}, P^{-}:=[Z]_{-} \cap[P]_{+}$. The restriction of a line $l$ or circle $c$ to $M_{Z}$ is denoted by $l^{*}:=\ln M_{Z}$ resp. $c^{*}:=c \cap M_{Z}$. For any two distinct points $P, Q \in M_{Z}$ we define

$$
P \sqcup Q:=\left\{\begin{array}{l}
l^{*} \quad \text { iff } P, Q \in l \in \mathscr{L}, \\
\{P\} \cup\left(P^{+}, P^{-}, Q\right)^{*} \quad \text { iff } P \text { and } Q \text { are nonparallel. }
\end{array}\right.
$$

Since two circles can have at most two points in common it follows that $P \bigsqcup Q=Q \sqcup P$ if and only if $P \sqcup Q=l^{*}$ for some $l \in \mathscr{L}$, provided the order $n$ of $\mathscr{M}$ is at least 5 . The verification of the axioms (L1), (L2) and (L3) (see [9]) is now straightforward. In order to define parallelism we have to require that the following condition holds in $\mathscr{M}$ for every point $Z$.

Let $P_{1}, Q_{1}, P_{2}, Q_{2} \in M_{z}$ and suppose that $P_{1}$ and $Q_{1}, P_{2}$ and $Q_{2}, P_{1}$ and $P_{2}$ are nonparallel. If there exists a circle $c$ touching $\left(P_{1}^{+}, P_{1}^{-}, Q_{1}\right)$ in $P_{1}^{-}$and touching ( $P_{2}^{+}, P_{2}^{-}, Q_{2}$ ) in $P_{2}^{+}$, then there also exists a circle $d$ touching ( $P_{1}^{+}, P_{1}^{-}, Q_{1}$ ) in $P_{1}^{+}$and touching ( $P_{2}^{+}, P_{2}^{-}, Q_{2}$ ) in $P_{2}^{-}$(see Figure 1).


Fig. 1.
In the definition of $P_{1} \sqcup Q_{1} \| P_{2} \sqcup Q_{2}$ we have to distinguish several cases.
Case 1. $P_{1}$ and $Q_{1}$ parallel, say $P_{1} \sqcup Q_{1}=l_{1_{1}}^{*}$ for some $l_{1} \in \mathscr{L}^{\varepsilon}$.

$$
P_{1} \sqcup Q_{1} \| P_{2} \sqcup Q_{2}: \Leftrightarrow P_{2} \sqcup Q_{2}=l_{2}^{*} \quad \text { for some } l_{2} \in \mathscr{L}^{e} .
$$

Case 2. $P_{1}$ and $Q_{1}$ nonparallel, $P_{1}, P_{2}$ parallel, say $P_{1}, P_{2} \in l \mathscr{L}^{\ell}$. From [9],

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proposition 3.1, it is clear that we have to define

$$
P_{1} \sqcup \dot{Q}_{1} \| P_{2} \sqcup Q_{2}: \Leftrightarrow P_{1} \sqcup Q_{1}=P_{2} \sqcup Q_{2}
$$

or

$$
\left(P_{1} \sqcup Q_{1}\right) \cap\left(P_{2} \sqcup Q_{2}\right)=\varnothing .
$$

Case 3. $P_{1}$ and $Q_{1}$ nonparallel and $P_{1}, P_{2}$ nonparallel. Put $P_{3}=\left[P_{1}\right]_{+} \cap$ $\cap\left[P_{2}\right]_{-}$and $P_{4}:=\left[P_{1}\right]_{-} \cap\left[P_{2}\right]_{+}$(see Figure 1). $P_{1} \sqcup Q_{1} \| P_{2} \sqcup Q_{2}: \Leftrightarrow$ $\Leftrightarrow$ There exists $P_{3} \sqcup Q_{3}$ such that

$$
\left(P_{3} \sqcup Q_{3}\right) \cap\left(P_{1} \cap Q_{1}\right)=\varnothing=\left(P_{3} \cap Q_{3}\right) \cap\left(P_{2} \cup Q_{2}\right) .
$$

Notice that condition ( $A$ ) is equivalent to: $P_{1} \sqcup Q_{1} \| P_{2} \sqcup Q_{2}$ implies $P_{2} \sqcup Q_{2} \| P_{1} \sqcup Q_{1}$, i.e., parallelism is a symmetric relation. We prove that parallelism is a transitive relation. Suppose $P_{1} L Q_{1} \| P_{2} \sqcup Q_{2}$ and $P_{2} \sqcup Q_{2} \|$ $P_{3} \sqcup Q_{3}$ (with distinct $P_{1}, P_{2}, P_{3}$ ). We prove that $P_{1} \sqcup Q_{1} \| P_{3} \sqcup Q_{3}$.
Case (a). $P_{1} \| Q_{1}$. Trivial
Case (b). $P_{1} \| Q_{1}, P_{2}, P_{3} \in l$ for some $l \in \mathscr{L}$. The transitivity follows at once from the following observation. If $c, d, e, \in \mathscr{C}$ and $c$ and $d$ touch in a point $P, d$ and $e$ touch in the same point $P$, then $c$ and $e$ touch in $P$. To show this suppose $Q \in c \cap e, Q \neq P$, then there are two circles through $Q$, namely $c$ and $e$, touching $d$ in $P$. This contradicts (M6).

Case (c). $P_{1} \| Q_{1}, P_{1} \in\left[P_{2}\right]_{\varepsilon}, P_{3} \in\left[P_{2}\right]_{-\varepsilon}$ for some $\varepsilon=+,-$. By definition $P_{1} \sqcup Q_{1} \| P_{3} \sqcup Q_{3}$.
Case (d). $P_{1}\left\|Q_{1}, P_{1}\right\|_{\varepsilon} P_{2}$ for some $\varepsilon=+,-, P_{3}\left\|P_{1}, P_{3}\right\| P_{2}$. Put $P_{4}:=$ $\left[P_{2}\right]_{\varepsilon} \cap\left[P_{3}\right]_{-\varepsilon} \cdot$ Since $P_{2} \sqcup Q_{2} \| P_{3} \sqcup Q_{3}$ there exists $Q_{4}$ such that $P_{2} \sqcup Q_{2} \|$ $P_{4} \sqcup Q_{4} \| P_{3} \sqcup Q_{3}$. Apply case (b) to find $P_{1} \sqcup Q_{1} \| P_{4} \sqcup Q_{4}$ and case (c) to find $P_{1} \sqcup Q_{1} \| P_{3} \sqcup Q_{3}$.
Case (e). $P_{1}\left\|Q_{1}, P_{1}\right\|_{\varepsilon} P_{3}$ for some $\varepsilon=+,-, P_{2}\left\|P_{1}, P_{2}\right\| P_{3}$. Put $P_{4}:=$ $=\left[P_{1}\right]_{\varepsilon} \cap\left[P_{2}\right]_{-\varepsilon}$. There exists $Q_{4}$ such that $P_{1} \sqcup Q_{1} \| P_{4} \sqcup Q_{4}$ and $P_{4} \sqcup Q_{4} \| P_{3} \sqcup Q_{3}$. Apply case (b).
Case (f). $P_{1} \| Q_{1}, P_{1}, P_{2}, P_{3}$ mutually nonparallel. Put $P_{4}:=\left[P_{1}\right]_{+} \cap\left[P_{2}\right]_{-}$. There exists $Q_{4}$ and that $P_{1} \sqcup Q_{1} \| P_{4} \sqcup Q_{4}$ and $P_{4} \sqcup Q_{4} \| P_{2} \sqcup Q_{2}$. Apply case (d) to find $P_{4} \sqcup Q_{4} \| P_{3} \sqcup Q_{3}$ and so $P_{1} \sqcup Q_{1} \| P_{3} \sqcup Q_{3}$.

Let $L^{Z}$ be the set of all $P \sqcup Q, P, Q \in M_{Z}, P \neq Q$. It is not hard to show that $\mathscr{M}^{Z}:=\left(M_{z}, L^{Z}, \amalg \|\right)$ satisfies all the axioms of a nearaffine plane except possibly ( $\mathbf{P} 2$ ) or ( $\mathbf{P} 2^{\prime}$ ). For ( $\mathbf{P} 2$ ) to hold we have to require:

Let $P_{1}, Q_{1}, P_{2}, Q_{2}$ be points as in (A). If $P_{1} \in\left(P_{2}^{+}, P_{2}^{-}, Q_{2}\right)$ and $P_{2} \in\left(P_{1}^{+}, P_{1}, Q_{1}\right)$. Then circles $c$ and $d$ as described in $(A)$ exist.
If we content ourself with the weaker ( $\mathrm{P} 2^{\prime}$ ) we have to require:
Let $\varepsilon$ be + or,$- A$ and $B$ to distinct points on $[Z]_{\varepsilon}, A \neq Z \neq B$

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and $c_{1}$ and $c_{2}$ two circles touching in A. Put (see Figure 2)

$$
\begin{aligned}
& C_{i}:=[Z]_{-\varepsilon} \cap c_{i}, \quad i=1,2, \\
& P_{i}:=[A]_{-\varepsilon} \cap\left[C_{i}\right]_{\varepsilon}, \quad i=1,2, \\
& Q_{i}:=[B]_{\varepsilon} \cap c_{i}, \quad i=1,2 \\
& D_{i}:=\left[Q_{i}\right]_{\varepsilon} \cap[Z]_{-\varepsilon}, \quad i=1,2, \\
& d_{i}:=\left(P_{i}, D_{i}, B\right), \quad i=1,2
\end{aligned}
$$

Then $d_{1}$ and $d_{2}$ touch in $B$.


Fig. 2.
If $\mathscr{M}$ is a Minkowski plane satisfying the conditions (A) and (B) or (A) and $(C)$ and $Z$ a point of $\mathscr{M}$, then the nearaffine plane $\mathscr{M}^{\mathbf{Z}}$ is called the residual plane with respect to $Z$.

For the remainder of the section let $\mathscr{M}=\left(M, \mathscr{L}^{+}, \mathscr{L}^{-}, \mathscr{C}\right)$ be a Minkowski plane satisfying the conditions (A) and (C). Since $\square$ and $\|$ are defined strictly in terms of the incidence in $\mathscr{M}$ it follows at once that an automorphism of $\mathscr{M}$ fixing a point $Z$, induces an automorphism of $\mathscr{M}^{Z}$, i.e., Aut $(\mathscr{M})_{Z} \lesssim$ Aut $\left(\mathscr{M}^{Z}\right)$. In fact, Aut $\left(\mathscr{M}_{\mathrm{Z}} \simeq\right.$ Aut $\left(\mathscr{M}^{Z}\right)$ as we shall see in a moment. The crucial observation is the following lemma.
3.1. LEMMA. Let $Z$ be a point of $\mathscr{M}$. For any two nonparallel points $A$ and $B$ of $M_{Z}$ let $[A, B]$ be the set of points consisting of $A, B, Z$ and the points $C \in M_{Z}$, nonparallel to $A$ and $B$,for which there is no set $P \sqcup Q \backslash\{P\}$ containing $A, B$ and C. Then

$$
[A, B]=(A, B, Z)
$$

Proof. Clearly both $[A, B]$ and $(A, B, Z)$ contain $A, B$, and $Z$. Let $C \in$ $(A, B, Z), C \neq A, B, Z$ then $(A, B, C)=(A, B, Z)$. Suppose for some $P, Q \in M_{Z}$ we have $A, B, C \in P \bigsqcup Q \backslash\{P\}$. Then $A, B, C \in\left(P^{+}, P^{-}, Q\right)\left\{P^{+}, P^{-}\right\}$, so $(A, B, C)=\left(P^{+}, P^{-}, C\right)$ a circle not passing through $Z$, a contradiction. Conversely, let $C \in[A, B], C \neq A, B, Z$ and suppose $C \in(A, B, Z)$. Then

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$Z \notin(A, B, C)$ and so $(A, B, C)$ intersects $[Z]_{+}$and $[\bar{Z}]_{-}$in points $P^{+}$and $P^{-}$ respectively, different from $Z$. So, with $P$ defined by $P=\left[P^{+}\right]_{-} \cap\left[P^{-}\right]_{+}$, $A, B, C$ are on $P \bigsqcup Q \backslash\{P\}$, a contradiction.

The lemma just proved shows that the residual plane $\mathscr{M}^{2}$ completely determines the Minkowski plane $\mathscr{M}$. The lines of $\mathscr{M}$ can be recovered from the straight lines of $\mathscr{M}^{Z}$, the circles not containing $Z$ from the proper lines of $M^{Z}$, and the circles containing $Z$ from the sets $[A, B]$. This proves the following theorem.

### 3.2. THEOREM. Let $Y$ and $Z$ be the points of $\mathscr{M}$. Then

(a) $\mathscr{M}^{Y} \simeq \mathscr{M}^{Z}$ iff there exists $\phi \in \operatorname{Aut}(\mathscr{M})$ such that $Y^{\phi}=Z$.
(b) Any automorphism of $\mathscr{M}^{Z}$ can be extended to an automorphism of $\mathscr{M}$ fixing $z$.
(c) Aut $\left(\mathscr{M}_{\mathcal{Z}} \simeq \operatorname{Aut}\left(\mathscr{M}^{2}\right)\right.$.

It is not hard to show that for any point $Z$ of $\mathscr{A}$ the residual plane $\mathscr{M}^{Z}$ satisfies the Veblen-condition $\left(V^{\prime}\right)$. In fact we can prove somewhat more.
3.3. THEOREM. Let $Z \in M, l \in \mathscr{L}, l \neq[Z]_{+},[Z]_{-}$and let $Y$ be defined by $Y=\ln \left([Z]_{+} \cup[Z]_{-}\right)$. Then $\mathscr{M}_{i^{*}}^{Z} \simeq \mathscr{M}_{Y}$, where $l^{*}$ is the straight line $\Lambda\{Y\}$ of $\mathscr{M}^{2}$ (notation as in [9]).

Proof. Define an isomorphism $\phi: M_{Z} \rightarrow M_{Y}$ of $\mathscr{M}_{Y}^{Z}$ onto $\mathscr{M}_{Y}$ as follows. For $P \in M_{z}, P \notin l^{*}$ we define $P^{\phi}:=P$, and for $P \in M_{z}, P \in l^{*}, P^{\phi}:=[P]_{-\varepsilon} \cap[Z]_{\varepsilon}$, where $\varepsilon$ is determined by $l \in \mathscr{L}^{\varepsilon}$.

As a direct consequence of this theorem we have the following result.
3.4. THEOREM. If the derived $\mathscr{M}_{Z}$ is a translation plane for every $Z \in M$, then the residual plane $\mathscr{M}^{2}$ is a nearaffine translation plane for every $Z \in M$.

Proof. Apply 3.3 and 5.2 of [9].
As a converse to this theorem we mention the following theorem.
3.5. THEOREM. Let $Z$ be a point of $\mathscr{M}$. if $\mathscr{M}^{z}$ is a nearaffine translation plane, then $\mathscr{M}_{z}$ is a translation plane and $\mathscr{M}^{z}$ and $\mathscr{M}_{z}$ have the same translation group.

Proof. By 3.2 every automorphism of $\mathscr{M}^{2}$ is also an automorphism of $\mathscr{M}_{Z}$, and it is not hard to show that a straight translation of $\mathscr{M}^{Z}$ with a direction corresponding to $\mathscr{L}^{E}$ is also a translation of $\mathscr{M}$. Let $\mathscr{T}_{+}$and $\mathscr{T}$. be the translation groups of $\mathscr{M}^{2}$ with directions $\mathscr{L}^{+}$and $\mathscr{L}^{-}$respectively. Since $\mathscr{F}_{+}$and $\mathscr{F}_{-}$are also translation groups of $\mathscr{M}_{Z}$ it follows that $\mathscr{T}_{+}$and $\mathscr{T}_{-}$ are elementary abelian. Hence, by 4.12 of of [ 9 ], the set $\mathscr{T}$ of all translation of $\mathscr{M}^{\mathrm{Z}}$ is a group and $\mathscr{T}=\mathscr{T}_{+} \mathscr{T}_{-}=$the full translation group of $M_{z}$.

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## 4. Characterizations of the known finite models

Using the correspondence with sharply triply transitive sets of permutations all known (finite) Minkowski planes can be described as follows. Let $P$ be a prime, $h$ a positive integer, $q:=p^{h}$ and $\phi$ an automorphism of $G F(q)$. Let $G(\phi)$ be the set of permutations acting on the projective line $\Omega:=P G(1, q)=G F(q) \cup$ $\cup\{\infty\}$ given by

$$
\begin{aligned}
& x \rightarrow \frac{a x+b}{c x+d}, \quad a, b, c, d \in G F(q), a d-b c=(\text { nonzero }) \text { square in } \\
& G F(q) \\
& x \rightarrow \frac{a x^{\phi}+b}{c x^{\phi}+d}, \quad a, b, c, d \in G F(q), a d-b c=\text { nonsquare in } G F(q),
\end{aligned}
$$

i.e., $G(\phi)=G_{1} \cup \phi G_{2}$, where $G_{1}:=\operatorname{PSL}(2, q)$ and $G_{2}:=\operatorname{PG}(2, q) \backslash \operatorname{PSL}(2, q)$. Then $G(\phi)$ is sharply triply transitive on $\Omega$ (cf. [7], [8], [10]). The residual planes of ( $\Omega, G(\phi)$ ) are easily seen to be the nearaffine translation planes described in [9], Section 8. We shall show that a Minkowski plane whose residual planes are nearaffine translation planes, is isomorphic to an $(\Omega, G(\phi))$.

Let $c$ be a circle of a Minkowski plane $\mathscr{M}$ of order $n$ and $Z$ a point of $\mathscr{M}, Z \notin c$. If $\mathscr{M}_{Z}$ is augmented to a projective plane, then the points of $c^{*}=$ $c \backslash\left([Z]_{+} \cup[Z]_{-}\right)$together with the two ideal points corresponding to $\mathscr{L}^{+}$and $\mathscr{L}^{-}$constitute an oval in this projective plane. In $n$ is even, there exists a point (the nucleus of the oval) in the projective plane such that the $n+1$ lines incident with this point are the $n+1$ tangents of the oval. If $n$ is odd, each point of the projective plane is incident with 0 or 2 tangents (see [3]). From this observation we deduce the following lemma.
4.1. LEMMA. Let $\mathscr{M}$ be a Minkowski plane of order $n$. If $n$ is even, there cannot exist 3 distinct circles $c_{1}, c_{2}, d$ such that $c_{1}$ and $c_{2}$ touch in a point $Z$ and $c_{i}$ touches $d$ in $P_{i} \neq Z, i=1,2$. In any case there cannot exist 4 distinct circles $c_{1}, c_{2}, c_{3}$ and $d$ such that $c_{1}, c_{2}, c_{3}$ touch in a point $Z$ and such that $c_{i}$ touches $d$ in a point $p_{i} \neq Z, i=1,2,3$.

Proof. Case $n$ is even. Suppose circles $c_{1}, c_{2}$ and $d$ as described exist. The lines $\left[[Z]_{+} \cap d\right]^{\ldots}$ and $\left[[Z]_{-} \cap d\right]_{+}$are tangents to the oval corresponding with $d$ in the projective plane associated with $\mathscr{A}_{\mathrm{z}}$. They intersect in a point of $M_{z}$. Also $c_{1}$ and $c_{2}$ are tangents to the oval. They intersect in an ideal point of the projective plane, a contradiction.

Case $n$ is odd. Now $c_{1}, c_{2}$ and $c_{3}$ correspond to tangents of the oval $d$ in the projective plane associated with $\mathscr{M}_{z}$. They intersect in one (ideal) point, a contradiction.
4.2. THEOREM. Let $\mathscr{M}=(\Omega, G)=\left(M, \mathscr{L}^{+}, \mathscr{L}^{-}, \mathscr{C}\right)$ be a Minkowski of order $n \geqslant 5$. Suppose conditions (A) and (C) hold in $\mathscr{M}$ and that $\mathscr{M}^{z}$ is a nearaffine translation plane for every point $Z$. Then $\mathscr{M} \simeq(\Omega, G(\phi))$.

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Proof. Fix $\alpha_{1} \in \Omega$. For each point $\left(\alpha_{1}, \beta\right) \in M$ there is an elementary Abelian group $\mathscr{T}_{\ldots}\left(\alpha_{1}, \beta\right)$ of translations of $\mathscr{A}^{\left(\alpha_{1}, \beta\right)}$ and $\mathscr{M}_{\left(\alpha_{1}, \beta\right)}$, and $\mathscr{T}_{-}\left(\alpha_{1}, \beta\right) \leqq$ Aut $(\mathscr{M})(3.2,3.4,3.5)$. Each $\mathscr{T}\left(\alpha_{1}, \beta\right)$ fixes all lines of $\mathscr{L}^{-}$and one lines of $\mathscr{L}^{+}$(namely the line $\{(\alpha, \beta) \mid \alpha \in \Omega\}$ ). Using the notation of Section 2, each $\mathscr{T}_{-}\left(\alpha_{1}, \beta\right)$ consists of positive automorphisms of the form $(1, b)$, where $b \in S^{\Omega}$ fixes $\beta$ and $G b=G$, i.e., for each $\beta \in \Omega$ there is an elementary Abelian group $B(\beta)$ which fixes $\beta$, acts regularly on $\Omega\{\beta\}$, and for which $G B(\beta)=G$. Define $B:=\langle B(\beta) \mid \beta \in \Omega\rangle$, then $B$ is doubly transitive on $\Omega$ and $G B=G$. Therefore, $G$ is a union of cosets of $B$ and in particular $B \subseteq G$. Hence, no nontrivial permutation in $B$ leaves 3 letters fixed. By a theorem of Feit ([4]), $B$ contains a normal subgroup of order $n+1$ or there exists an exactly triply transitive permutation group $B_{0}$ containing $B$ such that $\left[B_{0}: B\right] \leqslant 2$. Suppose $B$ contains a normal subgroup of order $n+1$, then $B$ also contains a sharply doubly transitive subgroup $B^{*}$. The circles $\left\{\left(\alpha, \alpha^{g}\right) \mid \alpha \in \Omega\right\}, g \in B^{*}$ together with the lines $l \in \mathscr{L}$ now constitute an affine plane of order $n+1$ and hence configuration as described in 4.1 exist, a contradiction. Therefore $B \leqslant B_{0}$, where $B_{0}$ is sharply 3-transitive, and $\left[B_{0}: B\right] \leqslant 2$. All sharply triply transitive groups are known (see [6]). If $n$ is even, then $B_{0} \simeq P S L(2, n)$ and so $B=G=$ $=\operatorname{PSL}(2, n)$, i.e. $\mathscr{A}$ is the classical Minkowski plane of order $n=2^{h}$. If $n$ is odd, there are at most two sharply 3-transitive groups of degree $n+1$ and such a group certainly contains $\operatorname{PSL}(2, n)$. The Sylow $p$-subgroups $B(\beta)$ of $B$ are the Sylow $p$-subgroups of $\operatorname{PSL}(2, n)$. Therefore $B \leqslant P S L(2, n)$ and since $|B| \geqslant$ $\geqslant \frac{1}{2}(n+1)(n)(n-1)$ it follows that $B \simeq \operatorname{PSL}(2, n)$. Thus, with $G_{1}:=\operatorname{PSL}(2, n)$ and $G_{2}:=P G L(2, n) \backslash P S L(2, n)$,

$$
G=G_{1} \cup \phi G_{2}
$$

for some $\phi \in S^{\Omega}$. It remains to show that $\phi$ is an automorphism of $G F(n)$. If $x, y$ and $z$ are three distinct points of $\Omega$, then there is a $g \in G_{1}$ such that $x^{\phi}=x^{g}$, $y^{\phi}=y^{g}, z^{\phi}=z^{g}$ for otherwise there exists $h \in G_{2}$ such that $x^{\phi}=x^{\phi h}$, $y^{\phi}=y^{\phi h}, z^{\phi}=z^{\phi h}$, i.e., $h=1$, contradicting $h \in G_{2}$. It follows that we may assume w.l.o.g. that $\phi$ fixes 0,1 and $\infty$. If we do so it also follows that

$$
\frac{x^{\phi}-y^{\phi}}{x-y}=\text { square in } G F(n) \text { for all } x, y \in G F(n), \quad x \neq y
$$

for $g \in G_{1}$ determined by $x^{\phi}=x^{\theta}, y^{\phi}=y^{g}, \infty=\infty^{\phi}=\infty^{g}$ has determinant $\left(x^{\phi}-y^{\phi}\right) /(x-y)$. By a theorem of Bruen and Levinger (see [2]) it follows that $\phi$ is an automorphism of $G F(n)$.

Using the previous theorems it is possible to give a geometric characterization of the Minkowski planes ( $\Omega, G(\phi)$ ). Consider the following configurational condition:
(D)

Let $\varepsilon$ be + or,$- l \in \mathscr{L}^{\varepsilon}$ and $V, W$ to distinct points on $l$. Suppose $c$ and $c^{\prime}$ are to distinct circles touching in $V$. Let $Y$ and $Q$ be two distinct points on $c, Y\|W, Q\| W$. Define

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$$
\begin{aligned}
& Y^{\prime}:=c^{\prime} \cap[Y]_{-\varepsilon^{\prime}} \\
& Q^{\prime}:=c^{\prime} \cap[Q]_{\varepsilon^{\prime}} \\
& d:=(Y, Q, W), \\
& d^{\prime}:=\left(Y^{\prime}, Q^{\prime}, W\right) .
\end{aligned}
$$

Then $d$ and $d^{\prime}$ touch in $W$ (see Figure 3).


Fig. 3.
Notice that (D) is nothing but a special case of the Desarques configuration (DI) in $\mathscr{M}^{Z}$ on the points $P, Q, R, P^{\prime}, Q^{\prime}, R^{\prime}$,
4.3. THEOREM. Let $\mathscr{M}$ be a Minkowski plane of order $n \geqslant 5$, and suppose (D) holds in $\mathscr{M}$. Then $\mathscr{M}$ is isomorphic to one of the planes ( $\Omega, G(\phi)$ ).

Of course the proof of 4.3 is based on 4.2 and it is clear that (D) implies (A). Also (C) is a consequence of (D).

### 4.4. LEMMA. Let $\mathscr{M}$ be a Minkowski plane of order $n$

(a) If $n$ is even then ( A ) implies ( B ) (hence (C)).
(b) In any case (D) implies (C).

Proof. (a) The following statement is easily seen to be equivalent to $(B)$ : If the circles $c$ and $d$ as described in (A) exist, then $P_{1} \in\left(P_{2}^{+}, P_{2}^{-}, Q_{2}\right) \Leftrightarrow$ $\Leftrightarrow P_{2} \in\left(P_{1}^{+}, P_{1}^{-}, Q_{1}\right)$. To prove this last statement, consider the configuration of condition (A) and suppose $c$ and $d$ exist, $P_{2} \in\left(P_{1}^{+}, P_{1}^{-}, Q\right)$ but $P_{1} \notin\left(P_{2}^{+}\right.$, $\left.P_{2}^{-}, Q_{2}\right)$. Let $e$ be the circle through $P_{1}$ touching $\left(P_{2}^{+}, P_{2}^{-}, Q_{2}\right)$ and $c$ in $P_{2}^{+}, f$ the circle through $P_{1}$ touching $\left(P_{1}^{+}, P_{1}^{-}, Q_{1}\right)$ in $P_{2}$. By (A) $e$ and $f$ touch in $P_{1}$. Similarly it follows that the circle $g$ through $P_{1}$ touching $\left(P_{2}^{+}, P_{2}^{-}, Q_{2}\right)$ in $P_{2}$ touches $f$ in $P_{1}$. Therefore, $g$ and $e$ touch in $P_{1}$ and so the circles $g, e$, $\left(P_{2}^{+}, P_{2}^{-}, Q_{2}\right)$ touch each other in $P_{2}^{+}, P_{2}^{-}, P_{1}$. This contradicts 4.1 since $n$ is even.
(b) Consider the configuration of condition (C). We claim that $\left(P_{1}, Q_{1}, Z\right)$ and ( $P_{2}, Q_{2}, Z$ ) touch in $Z$. If $\left(P_{i}, Q_{i}, Z\right)$ touches $c_{i}$ in $Q_{i}$ for $i=1,2$, this follows from (A). Suppose, therefore, that $\left(P_{1}, Q_{1}, Z\right)$ does not touch $c_{1}$ in $Q_{1}$, i.e., suppose that $\left(P_{1}, Q_{1}, Z\right)$ has another point $E_{1} \neq Q_{1}$ in common with

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$c_{1}$. Put $E_{2}:=\left[E_{1}\right]_{-\varepsilon} \cap c_{2}$. By (D) the circles $\left(E_{2}, Q_{2}, Z\right)$ and $\left(E_{1}, Q_{1}, Z\right)=$ $=\left(P_{1}, Q_{1}, Z\right)$ touch in $Z$. Suppose $\left(E_{2}, Q, Z\right)$ intersects $[A]_{-\varepsilon}$ in a point $P_{2}^{\prime} \neq P_{2}$ Let $Y$ be the point of intersection of $[Z]_{\varepsilon}$ and $\left(E_{2}, P_{2}^{\prime}, C_{2}\right)$. If we apply (D) twice it follows that $\left(E_{1}, P_{1}, Y\right)$ and ( $E_{1}, C_{1}, Y$ ) both touch $\left(E_{2}, P_{2}^{\prime}, C_{2}\right)$ in $Y$. Hence $\left(E_{1}, P_{1}, Y\right)=\left(E_{1}, C_{1}, Y\right)$ and impossibility because $P_{1} \| C_{1}$. We have proved $P_{2} \in\left(E_{2}, Q_{2}, Z\right)$, i.e., $\left(P_{1}, Q_{1}, Z\right)$ and $\left(P_{2}, Q_{2}, Z\right)$ touch in $Z$. So: $c_{1}$ and $c_{2}$ touch in $A$ implies $\left(P_{1}, Q_{1}, Z\right)$ and $\left(P_{2}, Q_{2}, Z\right)$ touch in $Z$. It is easily seen that the converse also holds. If we replace $c_{i}$ by $d_{i}, i=1,2$, it follows that $d_{1}$ and $d_{2}$ touch in $B$.

To finish the proof of 4.3 we have to show that all residual planes $\mathscr{M}^{2}$ are nearaffine transiation planes. By 3.4 it suffices to show that all derived planed $\mathscr{M}_{Z}$ are translation planes.
4.5. LEMMA. Let $\mathscr{M}$ be a Minkowski plane satisfying (D), then $\mathscr{M}_{z}$ is a translation plane for every point $Z$.

Proof. Let $Z \in M$ and $P, Q, R, P^{\prime}, Q^{\prime}, R^{\prime} \in M_{Z}$ such that $P\left\|_{-} P^{\prime}, Q\right\|_{-} Q^{\prime}$, $R \|_{-} R^{\prime}$, the line $P Q$ (in $\mathscr{M}_{z}$ ) is parallel to $P^{\prime} Q^{\prime}$ and $P R$ is parallel to $P^{\prime} R^{\prime}$. We have to show that $Q R$ is parallel to $Q^{\prime} R^{\prime}$, i.e., we have to show that the circles $(Z, Q, R)$ and $\left(Z, Q^{\prime}, R^{\prime}\right)$ touch in $Z$. We assume here that $P, Q, R$ (and also $P^{\prime}, Q^{\prime}, R^{\prime}$ ) are mutually nonparallel. The other cases follow from the cases we do consider. Put $Y=(P, Q, R) \cap[Z]_{+}$. If we apply (D) to $(P, Q, Z)$, $\left(P^{\prime}, Q^{\prime}, Z\right),(P, Q, Y)=(P, Q, R)$ and $\left(P^{\prime}, Q^{\prime}, Y\right)$, it follows that $(P, Q, R)$ and $\left(P^{\prime}, Q^{\prime}, Y\right)$ touch in $Y$. Application of $(\mathrm{D})$ to $(P, R, Z),\left(P^{\prime}, R^{\prime}, Z\right),(P, R, Y)=$ $=(P, Q, R)$ and $\left(P^{\prime}, R^{\prime}, Y\right)$ yields $(P, Q, R)$ and $\left(P^{\prime}, R^{\prime}, Y\right)$ touch in $Y$. Hence $\left(P^{\prime}, Q^{\prime}, Y\right)=\left(P^{\prime}, R^{\prime}, Y\right)=\left(P^{\prime}, Q^{\prime}, R^{\prime}\right)$. Finally we apply $(\mathrm{D})$ to $(Q, R, Y)\left(Q^{\prime}, R^{\prime}, Y\right)$, $(Q, R, Z)$ and $\left(Q^{\prime}, R^{\prime}, Z\right)$ and obtain thedesired result.

Notice that it is possible to give a proof of 4.3 without using the theory of nearaffine planes. Show directly, using (D), that any translation of a desired planes $\mathscr{M}_{Z}$ extends to an automorphism of $\mathscr{M}$. Then argue as we did in 4.2.

## REFERENCES

[^0]
## FINITE MINKOWSKI PLANES

7. Pedrini, C.: ${ }^{3}$-reti (non immergibli) aventi dei piani duali di quelli di Moulton quali sotto piani', Atti Accad. Naz. Lincei Rend. Sc. fis. mat. enat. 40 (1966). 385-392.
8. Quattrocchi. P.: 'Sugli insiemi di sostituzioni strettamente 3-transitivifiniti', Atti Sem. mat. fis. Univ. Modena 24(1975), 279-289 (1976).
9 Wilbrink, H. : 'Nearaffine Planes'. Geom. Dedicata 12 (1982) 53-62.
9. Wilbrink, H.: 'Nearaffine Planes and Minkowski Planes', Master's thesis, Tech Univ. Eindhoven, 1978.

## TWO-TRANSITIVE MINKOWSKI PLANES


#### Abstract

In this paper we determine all finite Minkowski planes with an automorphism group which satisfies the following transitivity property: any ordered pair of nonparallel points can be mapped onto any other ordered pair of nonparallel points.


## 1. Introduction

All known finite inversive planes have a two-transitive group of automorphisms. Conversely, every inversive plane admitting an automorphism group which is two-transitive on the points, is of a known type (cf. [9]).

For Minkowski planes the situation is quite similar. All known finite Minkowski planes have an automorphism group acting two-transitively on non-parallel points. In this note we shall show that this property is characteristic for the known Minkowski planes. More precisely, we shall prove the following theorem.

THEOREM. Let $\mathscr{M}$ be a finite Minkowski plane of odd order $n$, and suppose that $\mathscr{M}$ admits an automorphism group $\Gamma$ acting two-transitively on nonparallel points. Then $n$ is a prime power, $\mathscr{M} \simeq \mathscr{M}(n, \phi)$ for some field automorphism $\phi$ of $G F(n)$, and $\Gamma$ contains $\operatorname{PSL}(2, n) \times \operatorname{PSL}(2, n)$.

For a definition of $\mathscr{M}(n, \phi)$ see Section 2. As Minkowski planes of even order $n$ only exist for $n$ a power of 2 and are unique for given order $n=2^{a}$, this result completes the classification of the Minkowski planes with an automorphism group acting two-transitively on nonparallel points.

## 2. Definitions, notation and basic results

Let $M$ be a set of points and $\mathscr{L}^{+}, \mathscr{L}^{-}, \mathscr{C}$ three collections of subsets of $M$. The elements of $\mathscr{L}:=\mathscr{L}^{+} \cup \mathscr{L}^{-}$are called lines or generators, the elements of $\mathscr{C}$ are called circles. We say that $\mathscr{M}=\left(M, \mathscr{L}^{+}, \mathscr{L}^{-}, \mathscr{C}\right)$ is a $M$ inkowski plane if the following axioms are satisfied (cf. [8]).
(M1): $\quad \mathscr{L}^{+}$and $\mathscr{L}^{-}$are partitions of $M$.
(M2): $\quad\left|l^{+} \cap l^{-}\right|=1 \quad$ for all $l^{+} \in \mathscr{L}, l \in \mathscr{L}^{-}$.
(M3): Given any three points no two on a line, there is a unique circle passing through these three points.
(M4): $\quad|\cap c|=1$ for all $l \in \mathscr{L}, c \in \mathscr{C}$.
(M5) : There exist three points no two of which are on one line.

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(M6): Given a circle $c$, a point $P \in c$ and a point $Q \notin c, P$ and $Q$ not on one line, there is a unique circle $d$ such that $P, Q \in d$ and $c \cap d=$ $\{P\}$.

Two points $P$ and $Q$ are called plus-parallel (notation $P \|_{+} Q$ ) if $P$ and $Q$
 line of $\mathscr{L}^{-}$. Parallel (notation $P \| Q$ ) means either $P \|_{+} Q$ or $P \|_{-} Q$. For $P \in M$ we denote by $[P]_{+}$(resp. [P]_) the unique line in $\mathscr{L}^{+}$(resp. $\mathscr{L}^{-}$) incident with $P$. If $P, Q$ and $R$ are (distinct) nonparallel points, then we denote by $(P, Q, R)$ the unique circle containing $P, Q$ and $R$. Two circles $c$ and $d$ touch in a point $P$ if $c \cap d=\{P\}$.

We shall only consider finite Minkowski planes, i.e., Minkowski planes with a finite number of points. For finite Minkowski planes (M6) is a consequence of the other axiom (see [8]). It is easily seen that $\left|\mathscr{L}^{+}\right|=\left|\mathscr{L}^{-}\right|=$ $|l|=|c|=: n+1$ for all $l \in \mathscr{L}, c \in \mathscr{C}$. The integer $n$ is called the order of the Minkowski plane. Fix a point $P$ and put

$$
\begin{aligned}
M_{P} & : \\
L_{P} & \left.:=\left\{c^{*} \mid c \in \mathscr{C}, P \in c\right\} \cup\{ ]_{+} \cup[P]_{-}\right) \\
& \left.\mid \in \mathscr{L} \backslash\left\{[P]_{+},[P]_{-}\right\}\right\}
\end{aligned}
$$

where the * indicates that we have removed the point that the circle or line has in common with $[P]_{+} \cup[P]_{\ldots}$. Then $\mathscr{M}_{P}:=\left(M_{P}, L_{P}\right)$ is an affine plane with point set $M_{P}$ and line set $L_{P}$ (see, e.g., [8]). The projective plane associated with $\mathscr{M}_{P}$ will be denoted by $\mathscr{\mathscr { H }}_{p}$. We call $\mathscr{A}_{p}$ the derived plane with respect to the point $P$.

Following Benz [2], we sketch the close relationship between finite Minkowski planes and sharply triply transitive sets of permutations. Let $\Omega$ be a finite set, $|\Omega|=n+1$, and let $G$ be a subset of $\operatorname{Sym}(\Omega)$, the symmetric group on $\Omega$, acting sharply triply transitively on $\Omega$. Define

$$
\begin{aligned}
M & :=\Omega \times \Omega, \\
\mathscr{L}^{+} & :=\{\{(\alpha, \beta) \mid \alpha \in \Omega\} \mid \beta \in \Omega\}, \\
\mathscr{L}^{-} & :=\{\{(\alpha, \beta) \mid \beta \in \Omega\} \mid \alpha \in \Omega\}, \\
\mathscr{C} & :=\left\{\left\{\left(\alpha, \alpha^{g}\right) \mid \alpha \in \Omega\right\} \mid g \in G\right\} .
\end{aligned}
$$

Then $\mathscr{M}:=(\Omega, G):=\left(M, \mathscr{L}^{+}, \mathscr{L}^{-}, \mathscr{C}\right)$ is a Minkowski plane of order $n$. Conversely, every Minkowski plane can be obtained in this way.

Two Minkowski planes $\mathscr{M}=(\Omega, G)=\left(M, \mathscr{L}^{+}, \mathscr{L}^{-}, \mathscr{C}\right)$ and $\mathscr{M}^{\prime}=$ $\left(\Omega^{\prime}, G^{\prime}\right)=\left(M^{\prime}, \mathscr{L}^{+\prime}, \mathscr{L}^{-\prime}, \mathscr{C}^{\prime}\right)$ are said to be isomorphic if there is a bijection $s: M \rightarrow M^{\prime}$ such that

$$
\mathscr{L}^{s}=\mathscr{L}^{\prime} \quad \text { and } \quad \mathscr{C}^{s}=\mathscr{C}^{\prime}
$$

Either $\left(\mathscr{L}^{+}\right)^{s}=\mathscr{L}^{+\prime}$ and $\left(\mathscr{L}^{-}\right)^{s}=\mathscr{L}^{-\prime}$ or $\left(\mathscr{L}^{+}\right)^{s}=\mathscr{L}^{-\prime}$ and $\left(\mathscr{L}^{-}\right)^{s}=\mathscr{L}^{+\prime}$. In the first case, $s$ is called a positive isomorphism, in the second case, a negative isomorphism. If $s$ is a positive isomorphism then there exist bijections $a, b$ : $\Omega \rightarrow \boldsymbol{\Omega}^{\prime}$ such that $(\alpha, \beta)^{s}=\left(\alpha^{a}, \beta^{b}\right)$ for all $\alpha, \beta \in \Omega$, and $a^{-1} G b=G^{\prime}$. If $s$ is a negative isomorphism, then there exist bijections $a, b: \Omega \rightarrow \Omega^{\prime}$ such that $(\alpha, \beta)^{s}=\left(\beta^{b}, \alpha^{\alpha}\right)$ and $b^{-1} G^{-1} a=G^{\prime}$. It follows that we may assume w.l.o.g. that $G$ contains the identity permutation on $\Omega$.

A (positive, negative) automorphism of a Minkowski plane $\mathscr{M}$ is a (positive, negative) isomorphism of $A$ onto itself. The automorphism group

Aut $(\Omega, G) \leqslant \operatorname{Sym}(\Omega \times \Omega)$ of the Minkowski plane $(\Omega, G)$ is given by

$$
\begin{aligned}
\operatorname{Aut}(\Omega, G)= & \left\{(a, b) \in \operatorname{Sym}(\Omega) \times \operatorname{Sym}(\Omega) \mid a^{-1} G b=G\right\} \\
& \cup\left\{(a, b) \in \operatorname{Sym}(\Omega) \times \operatorname{Sym}(\Omega) \mid a^{-1} G b=G^{-1}\right\} \tau
\end{aligned}
$$

where $\tau \in \operatorname{Sym}(\Omega \times \Omega)$ is defined by $(\alpha, \beta)^{t}=(\beta, \alpha)$ for all $(\alpha, \beta) \in \Omega \times \Omega$.
We shall now describe all known finite Minkowski planes (cf. [14]).
Let $q$ be a prime power and let $\phi$ be a field automorphism of GF $(q)$. We shall denote by $\mu(q, \phi)$ the Minkowski plane $(\Omega, G)$ with $\Omega=P G(1, q)$, the projection line of order $q$, and $G$ the subset of Sym ( $\Omega$ ) consisting of the permutations

$$
x \mapsto \frac{a x+b}{c x+d}, \quad a d-b c=\text { a non-zero square of GF}(q)
$$

and

$$
x \mapsto \frac{a x^{\phi}+b}{c x^{\phi}+d}, \quad a d-b c=\text { a nonsquare of } \operatorname{GF}(q)
$$

Of course, if $q$ is even, we always get $G=\operatorname{PSL}(2, q)$, and it can be shown that these are the only Minkowski planes of even order (see [7]). For $q$ odd, $G$ is a group if and only if $\phi^{2}=1$ (see [10]), and nonisomorphic Minkowski planes of the same order $q$ can exist. Notice that $\mathscr{M}(q, \phi)$ has an automorphism group containing $\operatorname{PSL}(2, q) \times \operatorname{PSL}(2, q)$ which is two-transitive on nonparallel points, i.e., if $P, Q, P^{\prime}, Q^{\prime}$ are points such that $P \| Q$ and $P^{\prime} \| Q^{\prime}$, then there is an automorphism $g$ satisfying $P^{g}=P^{\prime}$ and $Q^{g}=Q^{\prime}$.

We conclude this section by listing some theorems on permutation groups which will be fundamental in our investigations. For the more standard results on (permutation) groups, the reader is referred to [11] or [17].
Result 1 (Gleason's lemma). Let $\Gamma$ be a permutation group of a finite set $M$ such that, for some prime $p$, every element of $M$ is fixed by a permutation in $\Gamma$ which has order $p$ and fixes no other element. Then $\Gamma$ is transitive on $M$ (see [5], 4.3.15, p.191).

A transitive permutation group which has the property that only the
identity fixes more than one letter, but the subgroup fixing a letter is nontrivial, is called a Frobenius group.

Result 2. In a Frobenius group the elements which fix no letter together with the identity form a regular normal subgroup (see [11], p. 495).

The regular normal subgroup in Result 2 is called the Frobenius kernel.
Result 3. Let $\Gamma$ be a 2-transitive permutation group on a finite set $M$ with an even number of letters such that only the identity fixes more than two letters. Then either $\Gamma$ contains a sharply 2 -transitive normal subgroup and $|M|$ is a power of 2 , or $\Gamma$ contains $\operatorname{PSL}(2, q)$ as a normal subgroup of index $\leqslant 2$ (see [6] and [12]).
Result 4. Let $\Gamma$ be a 2-transitive permutation group on a finite set $M$. If every element of $\Gamma$ which fixes an element of $M$ has odd order, then either $\Gamma$ is solvable (in which case $\Gamma$ is isomorphic to a subgroup of the group of semilinear transformations of a Galois field of characteristic 2 ) or $\Gamma$ contains as normal subgroup isomorphic to $\operatorname{PSL}(2, q)$ (see [1]).

## 3. Proof of theorem

For the proof of our theorem we require a number of lemmas. The first lemma shows that we can assume without loss of generality that an automorphism group which is two-transitive on nonparallel points, contains positive automorphisms only.
LEMMA 1. Let $\mathscr{M}=\left(M, \mathscr{L}^{+}, \mathscr{L}^{-}, \mathscr{C}\right)$ be a Minkowski plane of odd order $n$ and let $\Gamma^{*}$ be a group of automorphisms of $\mathscr{M}$ two-transitive on nonparallel points. Then $\Gamma:=\Gamma_{\mathscr{L}^{+}}=\Gamma_{\mathscr{L}^{-}}$is also two-transitive on nonparallel points ( $\Gamma_{\mathscr{L}^{+}}$is the set-wise stabilizer of $\mathscr{L}^{+}$in $\Gamma^{*}$ ).

Proof. Let $P$ and $Q$ be two points, $P \| Q$. Then

$$
\left[\Gamma_{P}: \Gamma_{P Q}\right]=\left[\Gamma_{P}^{*}: \Gamma_{P Q}^{*}\right]\left[\Gamma_{P Q}^{*}: \Gamma_{P Q}\right]\left[\Gamma_{P}^{*}: \Gamma_{P}\right]^{-1}
$$

Now $\quad\left[\Gamma_{P}^{*}: \Gamma_{P O}^{*}\right]=\left|M_{P}\right|=n^{2} \quad$ (as before $\quad M_{P}=M \backslash\left([P]_{+} \cup[P]_{-}\right)=$ $\{R \mid R \| P\}$ ), and $\left[\Gamma_{P Q}^{*}: \Gamma_{P Q}\right],\left[\Gamma_{P}^{*}: \Gamma_{P}\right] \in\{1,2\}$ since $\left[\Gamma^{*}: \Gamma\right] \in\{1,2\}$. Since $n$ is odd it follows that $\left[\Gamma_{P}: \Gamma_{P Q}\right]=n^{2}$, i.e., $\Gamma_{P}$ is transitive on $M_{P}$. Hence, $\Gamma$ is two-transitive on nonparallel points.

From now on $\mathscr{M}=\left(M, \mathscr{L}^{+}, \mathscr{L}^{-}, \mathscr{C}\right)=(\Omega, G)$ is a Minkowski plane of odd order $n \geqslant 5$ with a group $\Gamma$ of positive automorphisms acting twotransitively on nonparallel points. (For $n=3$ the theorem follows readily from [4].) We denote by $\Gamma\left(\mathscr{L}^{\varepsilon}\right)$ the subgroup of $\Gamma$ fixing all lines of $\mathscr{L}^{\varepsilon}$, $\varepsilon=+,-$. Notice that $\Gamma\left(\mathscr{L}^{-\varepsilon}\right)$ has a faithful representation on the $(n+1)$ lines of $\mathscr{L}^{\varepsilon}, \varepsilon=+,-$.

LEMMA 2. If $\Gamma\left(\mathscr{L}^{6}\right)$ contains $\operatorname{PSL}(2, n)$ for $\varepsilon=+$ or - , then $\mathscr{H} \simeq \mathscr{M}(n, \phi)$ for some $\phi \in A U T(\operatorname{GF}(n))$ and $\Gamma$ contains $\operatorname{PSL}(2, n) \times \operatorname{PSL}(2, n)$.

Proof. For convenience we take $\varepsilon=-1$. As a permutation group on $M=\Omega \times \Omega, \Gamma$ consists of permutations $\left(a_{\gamma}, b_{\gamma}\right) \in \operatorname{Sym}(\Omega) \times \operatorname{Sym}(\Omega)$ satisfying $a_{\gamma}^{-1} G b_{\gamma}=G, \gamma \in \Gamma$. Clearly, $\Sigma \leqslant \Gamma\left(\mathscr{L}^{-}\right)$is equivalent to $a_{\sigma}=1$ for all $\sigma \in \Sigma$. Hence $B:=\left\{b_{\sigma} \mid \sigma \in \Sigma\right\}$ is a subgroup of $\operatorname{Sym}(\Omega)$ satisfying $G B=G$. Therefore $G$ consists of a number of cosets of $B$, in particular $B \subseteq G$ since we are assuming that $1 \in G$. If $\Sigma \simeq B=G_{1}:=\operatorname{PSL}(2, n)$ then $G=G_{1} \cup \phi G_{2}$ for some $\phi \in \operatorname{Sym}(\Omega)$ where $G_{2}:=\operatorname{PGL}(2, n) \backslash G_{1}\left(\left|G_{1}\right|=\frac{1}{2}(n+1) n(n-1)\right.$ and $|G|=|\mathscr{G}|=(n+1) n(n-1)$ ). Viewing $\Omega$ as the projective line $G F(n) \cup\{\infty\}$ in the appropriate way, we claim that we may take $\phi \in \operatorname{Aut}(\mathrm{GF}(n))$. Let $x, y$ and $z$ be three distinct points of $\Omega$. Since $G$ is sharply triply transitive on $\Omega$, there exists a $g \in G$ such that $x^{\phi}=x^{g}, y^{\phi}=y^{g}$ and $z^{\phi}=z^{g}$. Suppose $g \in \phi G_{2}$, i.e., $g=\phi g_{2}$ for some $g_{2} \in G_{2}$, then $x^{\phi}=\left(x^{\phi}\right)^{g_{2}}$, $y^{\phi}=\left(y^{\phi}\right)^{g_{2}}, z^{\phi}=\left(z^{\phi}\right)^{q_{z}}$, and we get the contradiction $1=g_{2} \in G_{2}$.

We have shown: for any three distinct $x, y, z \in \Omega$ there is a $g_{1} \in G_{1}$ such that $x^{\phi}=x^{g_{1}}, y^{\phi}=y^{g_{1}}$ and $z^{\phi}=z^{g_{1}}$. It follows that we may assume without loss of generality that $\phi$ fixes 0,1 and $\infty$. If we do so it also follows that

$$
\frac{x^{\phi}-y^{\phi}}{x-y} \text { is a square in } \mathrm{GF}(n) \text { for all } x, y \in \mathrm{GF}(n), x \neq y
$$

for $g_{1} \in G_{1}$ determined by $x^{\phi}=x^{g}, y^{\phi}=y^{\theta}, \infty^{\phi}=\infty=\infty^{g}$ is the permutation $\left(\xi \mapsto\left(\left(x^{\phi}-y^{\phi}\right) /(x-y)\right)(\xi-y)+y^{\phi}\right) \in G_{1}$. By a theorem of Bruen and Levinger (see [3]) it follows that $\phi \in \operatorname{Aut}(\mathrm{GF}(n)$ ). It remains to show that $\Gamma\left(\mathscr{L}^{+}\right)$also contains $\operatorname{PSL}(2, n)$. Let $\gamma=\left(a_{\gamma}, b_{\gamma}\right) \in \Gamma$, then $a_{\gamma}^{-1} b_{\gamma} \in a_{\gamma}^{-1} G b_{\gamma}=$ $G \subseteq P \Gamma L(2, n)$. Hence,

$$
G_{1}^{a_{\nu}} \subseteq G^{a_{\gamma}}=a_{\gamma}^{-1} G a_{\gamma}=a_{\gamma}^{-1}\left(a_{\gamma} G b_{y}^{-1}\right) a_{\gamma}=G\left(a_{\gamma}^{-1} b_{\gamma}\right)^{-1} \subseteq \operatorname{PrL}(2, n) .
$$

Since $G_{1}^{a_{y}}$ is a two-transitive subgroup of $\operatorname{P\Gamma L}(2, n), G_{1}^{a_{\gamma}}$ contains $G_{1}$ so $G_{1}^{a_{\gamma}}=G_{1}$. Therefore $a_{\gamma} \in \operatorname{P\Gamma L}(2, n)$. Now $\left\{a_{\gamma} \mid \gamma \in \Gamma\right\}$ is a two-transitive subgroup of $\operatorname{P\Gamma L}(2, n)$, hence contains $G_{1}$. Since $a_{y}^{-1} b_{y} \in G=G_{1} \cup \phi G_{2}$ and $a_{\gamma}^{-1} G_{1} b_{\gamma}=G_{1}^{a_{\gamma}}\left(a_{\gamma}^{-1} b_{\gamma}\right)=G_{1}\left(a_{\gamma}^{-1} b_{\gamma}\right)$ either $a_{\gamma}^{-1} G_{1} b_{\gamma}=G_{1}$ or $a_{\gamma}^{-1} G_{1} b_{\gamma}=\phi G_{2}$, Since $G_{1}$ does not contain a subgroup of index $2,\left\{a_{\gamma} \mid \gamma \in \Gamma, a_{\gamma}^{-} G_{1} b_{\gamma}=G_{1}\right\}$ contains $G_{1}$. Let $a \in G_{1}$, then there is a $\gamma \in \Gamma$ such that $\gamma=(a, b), a^{-1} G_{1} b=$ $G_{1}$. Since $a \in G_{1}$ also $b \in G_{1}$. Hence $\left(1, b^{-1}\right) \in \Gamma$ and so $(a, 1)=(a, b)\left(1, b^{-1}\right) \in$ $\Gamma\left(\mathscr{L}^{+}\right)$.

LEMMA 3. Let $\varepsilon$ be + or -. If $\Sigma \leqslant \Gamma\left(\mathscr{L}^{\varepsilon}\right)$ is transitive on $\mathscr{L}^{-8}$ and $\Sigma_{l, m}=1$ for all $l, m \in \mathscr{L}^{-\varepsilon}, l \neq m$, then $\left|\Sigma_{l}\right| \leqslant 3$ for all $l \in \mathscr{L}^{-\varepsilon}$. If $\Gamma\left(\mathscr{L}^{\varepsilon}\right)_{l, m}=1$ for all $l, m \in \mathscr{L}^{-\varepsilon}$, then $\left|\Gamma\left(\mathscr{L}^{\varepsilon}\right)_{t}\right| \leqslant 3$ for all $l \in \mathscr{L}^{-\varepsilon}$.

Proof. Let $\Sigma \leqslant \Gamma\left(\mathscr{L}^{\theta}\right)$ be transitive on $\mathscr{L}^{-\varepsilon}$. Then $G$ contains a subgroup $H \simeq \Sigma$ (as permutation groups, see proof of Lemma 2). If $\Sigma_{t, m}=1$ for distinct $l, m$ in $\mathscr{L}^{-E}$, then $H_{\alpha, \beta}=1$ for distinct $\alpha, \beta \in \Omega$. It follows that the circles

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corresponding to the elements of $H$ cannot intersect each other in more than one point. Moreover, by Result 2, these circles fall into $\left|H_{\alpha}\right|$ classes of $n+1$ disjoint circles (each class corresponding to a coset of the Frobenius kernel of $H$ ). Thus each point is on $\left|H_{\alpha}\right|$ of these circles, one from each class, and circles in distinct classes intersect in exactly one point. Now, if $\left|\Sigma_{i}\right|=\left|H_{a}\right|>3$ we can find four circles $c_{1}, c_{2}, c_{3}$ and $d$ such that the $c_{i}$ touch each other in a point $P$ not on $d$ and such that the $c_{i}$ touch $d$ in three distinct points. However, this means that in the projective plane $\tilde{\mathscr{H}}_{p}$, the oval corresponding to $d$ has three tangents through a common point. As $n$, the other of $\tilde{\mathscr{M}}_{p}$, is odd, this is a contradiction.
Suppose $\Gamma\left(\mathscr{L}^{\ell}\right)_{\text {l, }}=1$ for distinct $l, m \in \mathscr{L}^{-\varepsilon}$. If $\Gamma\left(\mathscr{L}^{e}\right)=1$ for some (hence all) $l \in \mathscr{L}^{-\varepsilon}$ there is nothing to prove. If $\left|\Gamma\left(\mathscr{L}^{\ell}\right)\right|>1$, then $\Gamma\left(\mathscr{L}^{\ell}\right)$ is transitive on $\mathscr{L}^{-\varepsilon}$ by Result 1 and we can take $\Sigma=\Gamma\left(\mathscr{L}^{\varepsilon}\right)$.

LEMMA 4. Let $\varepsilon$ be + or - . If $\Gamma\left(\mathscr{L}^{\varepsilon}\right)$ is two-transitive on $\mathscr{L}^{-\varepsilon}$, then $n$ is a prime power, $\mathscr{A} \simeq \mathscr{M}(n, \phi)$ for some $\phi \in \operatorname{Aut}(\mathrm{GF}(n))$ and $\Gamma$ contains $\operatorname{PSL}(2, n \times$ PSL(2. $n$ ).

Proof. As $G$ is sharply triply transitive on $\Omega, \Gamma\left(\mathscr{L}^{\mathscr{f}}\right)_{l, m, n}=1$ for distinct lines $l, m, n \in \mathscr{L}^{-\varepsilon}$. By Result 3, either $\Gamma\left(\mathscr{L}^{e}\right)$ contains a sharply two-transitive subgroup, or $\Gamma\left(\mathscr{L}^{z}\right)$ contains $\operatorname{PSL}(2, n)$ as a normal subgroup of index $\leqslant 2$. The first alternative is impossible by Lemma 3. Lemma 2 now completes the proof.
LEMMA 5. If $\Gamma\left(\mathscr{L}^{\ell}\right)$ contains a nontrivial element fixing two lines of $\mathscr{L}^{-\varepsilon}$ $(\varepsilon=+$ or - ), then $n$ is a prime power, $\mathscr{M} \simeq \mathscr{M}(n, \phi)$ for some $\phi \in \operatorname{Aut}(\mathrm{GF}(n))$ and $\Gamma$ contains $\operatorname{PSL}(2, n) \times \operatorname{PSL}(2, n)$.

Proof. Suppose $1 \neq \gamma \in \Gamma\left(\mathscr{L}^{\varepsilon}\right)$ fixes $l, m \in \mathscr{L}^{-\varepsilon}, l \neq m$. We may assume that $\gamma$ has prime order. As remarked in the proof of Lemma 4, $\gamma$ fixes no other lines of $\mathscr{L}^{-8}$ besides $l$ and $m$. Since $\Gamma\left(\mathscr{L}^{\varepsilon}\right)$ is a normal subgroup of $\Gamma$, $\left\langle\gamma^{\alpha} \mid \alpha \in \Gamma_{1}\right\rangle \leqslant \Gamma\left(\mathscr{L}^{0}\right)$. By Result 1, it follows that $\left\langle\gamma^{\alpha} \mid \alpha \in \Gamma_{1}\right\rangle$ is transitive on $\mathscr{L}^{-\varepsilon} \backslash\{l\}$. Hence $\left\langle\gamma^{\alpha} \alpha \in \Gamma\right\rangle$ is two-transitive on $\mathscr{L}^{-\varepsilon}$. Now apply Lemma 4.

From the foregoing lemmas it is clear that our main objective will be to show that $\Gamma\left(\mathscr{L}^{e}\right)$ is nontrivial. For this it is necessary first to investigate how $\Gamma$ acts on $\mathscr{C}$ and how $\Gamma_{P}$ acts on $\mathscr{M}_{P}, P \in M$. Define a pencil to be any maximal set of mutually tangent circles through a common point $P$, called the carrier of the pencil. Thus the pencils with given carrier $P$ are essentially identical with parallel classes of lines in the affine plane $\mathscr{A}_{p}$. Every pencil contains $n$ circles. Every point is carrier of $n-1$ pencils.
LEMMA 6. For every point $P$ and pencil $\mathscr{P}$ with carrier $P, \Gamma_{P, \mathscr{g}}$ is transitive on the $n$ circles of $\mathscr{P}$.

Proof. Since $\Gamma$ is two-transitive on nonparallel points, $\Gamma_{P}$ is transitive on the points of $A_{p}$. By Theorem 3 of [16] we are done.

Thus, if circles $c$ and $d$ touch, then there exists $\gamma \in \Gamma$ such that $c^{\gamma}=d$. This shows that every $\Gamma$-orbit on $\mathscr{C}$ consists of à number of components of the touch-graph defined on $\mathscr{C}$ by:c, $d \in \mathscr{C}$ are adjacent iff $c$ and $d$ touch.

LEMMA 7. The touch-graph has 1 or 2 components. If it has 2 components, then each component contains $\frac{1}{2}(n+1) n(n-1)$ circles and every point is incident with $\frac{1}{2} n(n-1)$ circles of each component.

Proof. Let $c_{1}, c_{2}$ and $c_{3}$ be three distinct circles and $P$ a point, $P \notin c_{1}, c_{2}, c_{3}$. The ideal line of the affine plane $\mathscr{M}_{P}$ consists of the ideal points (i.e., parallel classes of $\left.\mathscr{M}_{P}\right) \mathscr{L}^{+} \backslash\left\{[P]_{+}\right\}, \mathscr{L}^{-} \backslash\left\{[P]_{-}\right\}$and the $(n-1)$ pencils with carrier $P$. The circles $c_{1}, c_{2}$ and $c_{3}$ correspond to ovals intersecting the ideal line in $\mathscr{L}^{+} \backslash\left\{[P]_{+}\right\}$and $\mathscr{L}^{-} \backslash\left\{[P]_{-}\right\}$. Thus, since $n$ is odd, for each $c_{i}$ there are $\frac{1}{2}(n-1)$ ideal points which are exterior with respect to $c_{i}$ (i.e., are the point of intersection of two tangents of $c_{i}$ ) and $\frac{1}{2}(n-1)$ ideal points which are interior with respect to $c_{i}$. This shows that at least two of $c_{1}, c_{2}$ and $c_{3}$ have an exterior point on the ideal line in common, hence are in the same component of the touch-graph. Therefore, the number of components is at most 2 . If there are 2 components and $c_{1}$ and $c_{2}$, say, are in distinct components, then the ideal points corresponding to the pencils fall into two classes: $\frac{1}{2}(n-1)$ are exterior with respect to $c_{1}$ and the other $\frac{1}{2}(n-1)$ are exterior with respect to $c_{2}$. Hence $P$ is incident with $\frac{1}{2} n(n-1)$ circles of each component and an easy counting argument shows that each component contains $\frac{1}{2}(n+1) n(n-1)$ circles.

Remark. The touch-graph of $\mathscr{M}(q, \phi), q$ odd, actually has two components.
By Lemmas 6 and 7, if $t$ is the number of $\Gamma$-orbits on $\mathscr{C}, t \in\{1,2\}$ and $\left[\Gamma: \Gamma_{c}\right]=t^{-1}(n+1) n(n-1)$ for all $c \in \mathscr{C}$. Using this result we can show the transitivity properties stated in the next lemma.

## LEMMA 8.

(i) If $c$ is a circle, then $\Gamma_{c}$ is two-transitive on $c$.
(ii) If $P$ is a point, then $\Gamma_{P}$ has $t$ orbits of length $t^{-1}(n-1)$ on the pencils with carrier $P$.
(iii) If $P$ and $Q$ are distinct points, $P \Downarrow Q$, then $\Gamma_{P, Q}$ has $t$ orbits of length $t^{-1}(n-1)$ on the circles containing $P$ and $Q$.
(iv) If $P$ and $Q$ are distinct points of the circle $c$, then $|\Gamma|=$ $=(n+1)^{2} n^{2}(n-1) t^{-1}\left|\Gamma_{P, Q .}\right|$.

Proof. Let $P$ and $Q$ be distinct points of the circle $c$, and let $\mathscr{P}$ be the pencil with carrier $P$ containing $c$. Denote by $s$ the number of pencils in the $\Gamma_{P}$-orbit containing $\mathscr{P}$. Then $\left[\Gamma_{P}: \Gamma_{P, \mathscr{g}}\right]=s$ and $\left[\Gamma_{P}: \Gamma_{P, c}\right]=n s$ by Lemma 6. Hence,

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$$
\begin{aligned}
(n+1) \geqslant\left[\Gamma_{c}: \Gamma_{c, p}\right] & =\frac{\left|\Gamma_{c}\right|}{|\Gamma|} \left\lvert\, \frac{|\Gamma|}{\left|\Gamma_{P}\right|} \cdot \frac{\left|\Gamma_{P}\right|}{\left|\Gamma_{P, t}\right|}\right. \\
& =\frac{1}{t^{-1}(n+1) n(n-1)} \cdot(n+1)^{2} \cdot n s \\
& =\frac{s t(n+1)}{n-1}=s t+\frac{2 s t}{n-1} .
\end{aligned}
$$

Thus, $s t=\frac{1}{2}(n-1) u$ with $u \in \mathbb{N}$, and so $(n+1) \geqslant\left[\Gamma_{c}: \Gamma_{c, p}\right]=\frac{1}{2}(n+1) u$, i.e., $u \in\{1,2\}$. As $s=\frac{1}{2} t^{-1}(n-1) u$ with $u, t \in\{1,2\}$ and $n$ is odd, $(n, s)=1$. Therefore it follows from

$$
\begin{aligned}
n & \geqslant \frac{\left|\Gamma_{c, P}\right|}{\left|\Gamma_{c, P, Q}\right|}=\frac{\left|\Gamma_{P, c}\right|}{\left|\Gamma_{P}\right|} \cdot \frac{\left|\Gamma_{P}\right|}{\left|\Gamma_{P, Q}\right|} \cdot \frac{\left|\Gamma_{P, Q}\right|}{\left|\Gamma_{P, Q, Q}\right|} \\
& =\frac{1}{n s} \cdot n^{2} \cdot\left[\Gamma_{P, Q}: \Gamma_{P, Q, c}\right]=\frac{n}{s}\left[\Gamma_{P, Q}: \Gamma_{P, Q, c}\right]
\end{aligned}
$$

that $\left[\Gamma_{c, P}: \Gamma_{c, P, Q}\right]=n$ and $\left[\Gamma_{P, Q}: \Gamma_{P, Q, c}\right]=s$. Now from $\left[\Gamma_{c, P}: \Gamma_{c, P, Q}\right]=n$ it follows that $\Gamma_{c, P}$ is transitive on $c \backslash\{P\}$, hence, since $P$ was an arbitrary point of $c, \Gamma_{c}$ is two-transitive on $c$. Therefore $(n+1)=\left[\Gamma_{c}: \Gamma_{c, p}\right]=\frac{1}{2}(n+1) u$, so $u=2$ and $s=t^{-1}(n-1)$. Finally,

$$
|\Gamma|=\frac{|\Gamma|}{\left|\Gamma_{P}\right|} \cdot \frac{\left|\Gamma_{P}\right|}{\left|\Gamma_{P, Q}\right|} \cdot \frac{\left|\Gamma_{P, Q}\right|}{\left|\Gamma_{P, Q, c}\right|} \cdot\left|\Gamma_{P, Q, c}\right|=(n+1)^{2} n^{2}(n-1) t^{-1} \cdot\left|\Gamma_{P, Q, c}\right|
$$

which proves (iv).
LEMMA 9. Let $P$ be a point. If $\Gamma_{p}$ has odd order, then $n$ is a power of a prime, $\mathscr{M} \simeq \mathscr{M}(n, \phi)$ for some $\phi \in \operatorname{Aut}(\mathrm{GF}(n))$ and $\Gamma$ contains $\operatorname{PSL}(2, n) \times \operatorname{PSL}(2, n)$.

Proof. Fix a line $l \in \mathscr{L}^{+}$and let $\Delta \simeq \Gamma_{i} /\left(\Gamma\left(\mathscr{L}^{-}\right) \cap \Gamma_{t}\right)$ be the permutation group on $l$ induced by $\Gamma_{l}$. As $\Gamma$ is two-transitive on the nonparallel points of $\mathscr{A}, \Delta$ is two-transitive on $l$. As $\Gamma_{P}$ has odd order, $\Delta_{p}$ has odd order for all $P \in l$. By Result 4 , either $\Delta$ is solvable or $\Delta$ contains $\operatorname{PSL}(2, n)$ as a normal subgroup. If $\Delta$ is solvable, then $\Delta$ is isomorphic to a subgroup of the group of semilinear transformations of a Galois field of characteristic 2, i.e., $n+1=$ $2^{a}$ for some $a \in \mathbb{N}$ and $\mid \Delta \|(n+1) n a$. If $\Delta$ contains $\operatorname{PSL}(2, n)$ as a normal subgroup, then $n=p^{b}$ for some prime $p$ and $b \in \mathbb{N}$ and $\Delta$ is a subgroup of $\operatorname{P\Gamma L}(2, n)$, i.e., $|\Delta| \mid(n+1) n(n-1) b$. By Lemma $8(i v)$, the order of $\Gamma_{i}$ is $(n+1) n^{2}(n-1) t^{-1} \cdot\left|\Gamma_{P, Q, c}\right|$. In both cases it follows from $n \geqslant 5$ that $\mid \Gamma\left(\mathscr{L}^{-}\right) \cap$ $\Gamma_{i}\left|=\left|\Gamma\left(\mathscr{L}^{-}\right)_{i}\right|>3\right.$.

By Lemma 3 there exists a nontrivial element of $\Gamma\left(\mathscr{L}^{-}\right)$fixing two distinct lines of $\mathscr{L}^{+}$. Lemma 5 now completes the proof.

By the previous lemma we may assume from now on that $\Gamma_{P}$ has even
order. More in particular, $\Gamma_{P}$ contains involutions. Since $n$ is odd, every involution $\tau \in \Gamma_{P}$ either induces a homology of the projective plane $\mathscr{M}_{P}$ associated with the affine plane $\mathscr{M}_{p}$, or the $\tau$-fixed points and lines of $\tilde{\mathscr{M}}_{P}$ constitute a Baer subplane of $\tilde{\mathscr{M}}_{P}$ (cf. [5], p. 172). Our next lemma deals with the case where $\Gamma_{P}$ contains a homology.

LEMMA 10. Let $P \in M$ and suppose that $\tau \in \Gamma_{P}$ is an involution which, considered as a collineation of $\tilde{\mathscr{M}}_{P}$, is a homology. Then $n$ is a prime power, $\mathscr{M} \simeq \mathscr{M}(n, \phi)$ for some $\phi \in \operatorname{Aut}(\mathrm{GF}(n))$ and $\Gamma$ contains $\operatorname{PSL}(2, n) \times \operatorname{PSL}(2, n)$. If $\Gamma_{P}$ has even order and
(i) $n$ is not a square, or
(ii) $\quad t=1$ (i.e., $\Gamma$ is transitive on $\mathscr{C}$ ), then $\Gamma_{P}$ contains homologies.

Proof. We distinguish two cases:
Case (a). The axis of $\tau$ is the ideal line of $\mathscr{M}_{P}$. Now, since $\Gamma_{P}$ is transitive on $M_{P}, \mathscr{M}_{P}$ is a translation plane and $\Gamma_{P}$ contains the full translation group of $\mathscr{M}_{P}$ (see [5], p. 187, result 4.3.1). Let $\Sigma^{(P)}$ be the subgroup of $\Gamma_{P}$ consisting of those translations of $\mathscr{M}_{P}$ which fix all lines of $\mathscr{L}^{-}$. Then $\Sigma^{(P)}$ is transitive on $\mathscr{L}^{+} \backslash\left\{[P]_{+}\right\}$, hence $\Sigma:=\left\langle\Sigma^{(P)} \mid P \in M\right\rangle$ is two-transitive on $\mathscr{L}^{+}$. Since $\Sigma \leqslant \Gamma\left(\mathscr{L}^{-}\right)$we are done by Lemma 4.

Case (b). The axis of $\tau$ is an affine line of $\mathscr{M}_{\mathbf{P}}$. Clearly, the axis of $\tau$ corresponds to aline $l \neq[P]_{+},[P]_{-}$of $\mathscr{M}$, say $l \in \mathscr{L}^{+} \backslash\left\{[P]_{+}\right\}$. Now $1 \neq \tau \in \Gamma\left(\mathscr{L}^{-}\right)$ and $\tau$ fixes the two distinct lines $[P]_{+}$and $l$ of $\mathscr{L}^{+}$. By Lemma 5 we have completed the proof of our first claim.

The order of a Baer subplane of $\mathscr{M}_{P}$ is $\sqrt{n}$. Hence, if $n$ is not a square, every involution in $\Gamma_{P}$ acts as a homology of $\tilde{\mathscr{M}}_{P}$. Suppose $t=1$. Let $\Lambda$ be a Sylow 2-subgroup of $\Gamma_{P}$ and let $\tau$ be an involution in the center of $\Lambda$. Suppose the $\tau$-fixed points and lines of $\tilde{\mathscr{M}}_{P}$ constitute a Baer subplane. The two ideal points of $\mathscr{M}_{P}$ corresponding to $\mathscr{L}^{+}$and $\mathscr{L}^{-}$are fixed by $\Gamma_{P}$, and by Lemma 8 (ii) $\Gamma_{P}$ is transitive on the remaining $n-1$ ideal points. Let $2^{a} \|(n-1)$. By [17], Theorem 3.4, every $\Lambda$-orbit on these $n-1$ ideal points has length divisible by $2^{a}$. The ideal line of $\mathscr{M}_{p}$ is fixed by $\tau$ and contains therefore, apart from the ideal points corresponding to $\mathscr{L}^{+}$and $\mathscr{L}^{-}, \sqrt{n}-1$ fixed points. Since $\tau \in Z(\Lambda), \Lambda$ permutes these $\sqrt{n}-1$ points. However, $2^{b} \|(\sqrt{n}-1)$ with $b<a$, contradicting the fact that each of these $\sqrt{n}-1$ points is in a $\Lambda$-orbit with length divisible by $2^{a}$.

For the proof of our main result we need one more definition and lemma.
DEFINITION. Suppose $M_{1} \subseteq M ; \mathscr{L}_{1}^{\varepsilon} \subseteq \mathscr{L}^{\varepsilon}, \varepsilon=+,-; \mathscr{C}_{1} \subseteq \mathscr{C}$. Put $\mathscr{L}_{1}^{\varepsilon *}:=\left\{l \cap M_{1} \mid l \in \mathscr{L}_{1}^{\varepsilon}\right\}, \quad \varepsilon=+,-; \mathscr{C}_{1}^{*}: \doteq\left\{c \cap M_{1} \mid c \in \mathscr{C}_{1}\right\}$. If $\mathscr{M}_{1}:=$ $\left(M_{1}, \mathscr{L}_{1}^{+*}, \mathscr{L}_{1}^{-*}, \mathscr{C}_{1}^{*}\right)$ is a Minkowski plane with the property that any

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two circles which touch in $\mathscr{M}_{1}$, touch in $\mathscr{M}$, then $\mathscr{M}_{1}$ is called a subplane of $\mathscr{A}$ (compare [5], p. 258).
LEMMA 11. Let $\Delta$ be a group of positive automorphisms of $\mathscr{M}$. Let $M_{1}$ be the set of points left fixed by $\Delta ; \mathscr{L}_{1}^{+}$(resp. $\left.\mathscr{L}^{-}\right)$the set of lines of $\mathscr{L}^{+}\left(\right.$resp. $\left.\mathscr{L}^{-}\right)$ left fixed by $\Delta$; and $\mathscr{C}_{1}$ the set of circles left fixed by $\Delta$. Then $\mathscr{A}_{1}:=$ $\left(M_{1}, \mathscr{L}_{1}^{+*}, \mathscr{L}_{1}^{-*}, \mathscr{C}_{1}^{*}\right)$ is a subplane of $\operatorname{MI}$ if and only if $M_{1}$ contains, at least three mutually nonparallel points.

Proof. Straightforward verification.
We are now ready to prove our main result.
THEOREM. Let $\mathscr{A}=\left(M, \mathscr{L}^{+}, \mathscr{L}^{-}, \mathscr{C}\right)$ be a finite Minkowski plane of odd order $n$, and suppose that, $\mathscr{M}$ admits an automorphism group $\Gamma$ two-transitive on nonparallel points. Then $n$ is a prime power, $\mathscr{M} \simeq \mathscr{M}(n, \phi)$ for some $\phi \in$ $\operatorname{Aut}(\mathrm{GF}(n))$ and $\Gamma$ contains $\operatorname{PSL}(2, n) \times \operatorname{PSL}(2, n)$.

Proof. Suppose $\mathscr{A}$ is a counter example to the theorem of minimal order. By Lemma 1 we may assume that $\Gamma$ contains positive automorphisms only. By Lemma $9, \Gamma_{P}$ has even order for all $P \in M$. By Lemma 10 every involution in $\Gamma_{P}$ has $(\sqrt{n}+1)^{2}$ fixed points. Hence, if $\Lambda$ is a 2 -subgroup of $\Gamma$ maximal with respect to fixing at least three mutually nonparallel points, $\Lambda \neq 1$. Let $\mathscr{M}_{1}=\left(M_{1}, \mathscr{L}_{1}^{+*}, \mathscr{L}_{1}^{-*}, \mathscr{B}_{1}^{*}\right)$ be the subplane of $\mathscr{M}$ consisting of the $\Lambda$-fixed points, lines and circles of $\mathscr{M}$ of order $n_{1}$, say. Clearly $n_{1}$ is odd, and since $\Lambda \neq 1$ we have $n_{1}<n$. We claim that $N_{\Gamma}(\Lambda)$, considered as an automorphism group of $\mathscr{M}_{1}$, acts two-transitively on the nonparallel points of $\mathscr{M}_{1}$. To see this, let $c \in \mathscr{C}_{1}$. Then $\Lambda \leqslant \Gamma_{c}$ and $\Lambda$, considered as a permutation group on $c$, is a 2 -subgroup of $\Gamma_{c}$ maximal with respect to fixing at least three points of $c$. By Lemma $8(\mathrm{i}), \Gamma_{c}$ is two-transitive on $c$, hence $N_{\Gamma_{c}}(\Lambda)$ is two-transitive on $c^{*}:=c \cap M_{1}$ (see [1], Lemma 3.3). Now let $A_{1}, A_{2}$ and $B_{1}, B_{2}$ be two pairs of nonparallel points of $\mathscr{M}_{1}$. If $A_{i} \| B_{j}, i, j=1,2$, and $c_{1}$ is the unique circle containing $A_{2}, B_{1}, B_{2}$, and $c_{2}$ is the unique circle containing $A_{2}, B_{1}, B_{2}$, then there is a $\gamma_{1} \in N_{\Gamma_{e_{e}}}(\Lambda)$ and a $\gamma_{2} \in N_{\Gamma_{c_{2}}}(\Lambda)$ such that $\boldsymbol{A}_{1}^{\gamma_{1}}=\boldsymbol{A}_{2}, \boldsymbol{A}_{2}^{\gamma_{1}}=\boldsymbol{B}_{1}, \boldsymbol{A}_{1}^{\gamma_{2}}=\boldsymbol{B}_{1}, B_{1}^{\gamma_{2}}=\boldsymbol{B}_{2}$. Hence $\gamma=\gamma_{1} \gamma_{2} \in N_{\mathrm{r}}(\Lambda)$ satisfies $A_{1}^{\gamma}=B_{1}$ and $A_{2}^{\gamma}=B_{2}$. Repeated application of this result in case $A_{i} \| B_{j}$ for some $i$ and $j$, proves our claim. Since $\mathscr{M}$ was supposed to be a minimal counter example, $n_{1}$ is a prime power, say $n_{1}=p^{a}$ with $p$ prime and $a \in \mathbb{N}$. If $P \in M_{1}$, then the projective plane $\left(\tilde{\mathscr{H}}_{1}\right)_{P}$ associated with $\left(\mathscr{M}_{1}\right)_{P}$ is a subplane of the projective plane $\tilde{\mathscr{M}}_{p}$ associated with $\mathscr{M}_{p}$ (this is why we required in the definition of a subplane of a Minkowski plane, that circles tangent in $\mathscr{M}_{1}$ are also tangent in $\mathscr{M})$. In fact $\left(\tilde{\mathscr{M}}_{1}\right)_{p}$ is a 2 -subplane of $\mathscr{M}_{p}$ is the sense of Ostrom and Wagner [15]. By their Theorem $6, n=n_{1}^{2 g}$ for some integer $g$. Hence, also $n$ is a prime power, $n=p^{b}$ with $b=a 2^{g}$. Let $\Pi$ be a Sylow $p$-subgroup of $\Gamma_{P}, P \in M$. Let $\pi$ be an element in the centre of $\Pi$. Since $\pi$ fixed the two ideal points corresponding to $\mathscr{L}^{+}$and $\mathscr{L}^{-}$of $\mathscr{H}_{p}, \pi$ also fixes an affine

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line $L$ of $\mathscr{M}_{p}$. Suppose $L$ intersects the ideal line of $\mathscr{M}_{p}$ in a point $A$. Then $\Pi$ fixes $A$ for if $A^{\sigma} \neq A$ for some $\sigma \in \Pi$, then $L^{\sigma}$ and $L$ intersect in an affine point $Q$ of $M_{p}$. Since $\Pi$ permutes the fixed objects of $\pi, L^{\sigma}$ hence $Q$ is fixed by $\pi$. Since $\Gamma_{P}$, hence $\Pi$, is transitive on the $n^{2}$ affine points of $\mathscr{M}_{P}$, every affine point of $\mathscr{M}_{P}$ is fixed by $\pi$, i.e., $\pi=1$ a contradiction. By Theorem 3 of [16] $\Gamma_{P, A}$, hence $\Pi$ is transitive on the $n$ affine lines through $A$. Therefore, $\pi$ fixes all lines through $A$, i.e., $\pi$ is an elation of $\tilde{\mathscr{H}}_{P}$ with centre $A$ and axis the ideal line of $\mathscr{A}_{\mathrm{P}}$. Suppose $A$ is the ideal point corresponding to $\mathscr{L}^{-\varepsilon}$ for $\varepsilon=+$ or - , then $\pi \in \Gamma\left(\mathscr{L}^{-\varepsilon}\right)_{{ }_{P P_{c}}}$. By Lemma $5, \Gamma\left(\mathscr{L}^{-\varepsilon}\right)_{l_{, m}}=1$ for distinct lines $l, m \in \mathscr{L}^{\varepsilon}$, so by Lemma 3, $p \leqslant$ order of $\pi \leqslant\left|\Gamma\left(\mathscr{L}^{-\varepsilon}\right)_{P_{p l}}\right| \leqslant 3$, i.e., $p=3$. Also $\Gamma\left(\mathscr{L}^{-\varepsilon}\right)$ is a Frobenuis group on the $(n+1)$ lines of $\mathscr{L}^{E}, \Gamma\left(\mathscr{L}^{-\varepsilon}\right) \leqslant \Gamma$ and $\Gamma$ acts twotransitively on $\mathscr{L}^{t}$, hence the Frobenius kernel of $\Gamma\left(\mathscr{L}^{-\varepsilon}\right)$ is an elementary abelian 2-group and in particular $n+1=2^{c}$ for some $c \in \mathbb{N}$. However, $n+1=$ $p^{b}+1=3^{a 2^{g}}+1 \equiv 2(4)$ and so we have shown that $A$ is an ideal point corresponding to a pencil with carrier $P$. Let $T$ be the group of translations of $\mathscr{M}_{P}$ contained in $\Gamma_{p}$ and for each pencil $\mathscr{P}$ with carrier $P$ let $T(\mathscr{P})$ be the group of translations of $T$ fixing all circles of $\mathscr{P}$. By Lemma 10 and Lemma 8 (ii), $\Gamma_{p}$ has two orbits of length $\frac{1}{2}(n-1)$ on the pencils with carrier $P$. Put $x=|T(\mathscr{P})|$ for $\mathscr{P}$ in the first, and $y=|T(\mathscr{P})|$ for $\mathscr{P}$ in the second orbit. It follows that

$$
\begin{align*}
|T| & =1+(x-1) \cdot \frac{1}{2}(n-1)+(y-1) \cdot \frac{1}{2}(n-1)=  \tag{1}\\
& =1+\frac{1}{2}(x+y-2)(n-1),
\end{align*}
$$

and one of $x$ and $y \geqslant p$, so $x+y \geqslant p+1$. Also, if $s$ is the number of $T$-orbits on $M_{p}$,

$$
\begin{equation*}
s|T|=n^{2} . \tag{2}
\end{equation*}
$$

Since $x+y \geqslant p+1 \geqslant 4$, it follows that $|T| \geqslant n$, hence $s \leqslant n$. From (1) and (2) it also follows that $s \equiv 1\left(\bmod \frac{1}{2}(n-1)\right)$. Since $T$ is not transitive on $M_{P}$, $s>1$. Therefore $s=n,|T|=n$ and $p=3$. We list some properties of $T$.
(i) As a translation group containing translations in different directions, $T$ is elementary abelian,
(ii) $T \triangleleft \Gamma_{p}$,
(iii) $T$ acts regularly on the lines of $\mathscr{L} \succeq\left\{[P]_{\varepsilon}\right\}, \varepsilon=+,-$,
(iv) the subgroups $\langle\tau\rangle, \tau \in T$ are in $1-1$ correspondence with the $\frac{1}{2}(n-1$ ) pencils with carrier $P$ in a $\Gamma_{P}$-orbit: $\tau \leftrightarrow$ pencil $\mathscr{P}$ iff centre of $\tau=\mathscr{P}$; $\Gamma_{P}$ acts on this orbit as $\Gamma_{P}$ acts on $\{\langle\tau\rangle \mid \tau \in T\}$ by conjugation.
Take $Q \in M_{P}$. By Lemma 8(iii), $\Gamma_{P, Q}$ is still transitive on the pencils with carrier $P$ in a $\Gamma_{P}$-orbit, so $\Gamma_{P, Q}$ acts by conjugation transitively on the subgroups $\langle\tau\rangle, \tau \in T$. By (ii) and (iii), $T$ is a regular normal subgroup of $\Gamma_{P}$ considered as a permutation group on $\mathscr{L}^{+} \backslash\left\{[P]_{+}\right\}$. Since $\Gamma_{P, Q} \leqslant \Gamma_{P,[\uparrow] \mathrm{c}}$,

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$\Gamma_{P, Q}$ acts on $\mathscr{L}^{+} \backslash\left\{[P]_{+},[Q]_{+}\right\}$as it does on $T \backslash\{1\}$ by conjugation. It follows that either $\Gamma_{P, Q}$ is transitive or has two orbits of length $\frac{1}{2}(n-1)$ on $\mathscr{L}^{+} \backslash\left\{[P]_{+},[Q]_{+}\right\}$. The former alternative is impossible: an involution in the center of a Sylow 2-subgroup of $\Gamma_{P}$ is a homology (see the last part of the proof of Lemma 10). Therefore $\Gamma_{P, Q}$ has 2 orbits of length $\frac{1}{2}(n-1)$ on $\left.\mathscr{L}^{+} \backslash[P]_{+},[Q]_{+}\right\}$and it acts on both orbits as it acts on the subgroups $\langle\tau\rangle, \tau \in T$ by conjugation. Let $c$ be a circle through $P$ and $Q$ in the pencil $\mathscr{P}$, where $\mathscr{P}$ is the centre of $\langle\tau\rangle$, say. Then $\Gamma_{P, Q, c}$ fixes $\mathscr{P}$, hence $\Gamma_{P, Q, c}$ fixes $\langle\tau\rangle$ by conjugation and therefore also two distinct lines $l, m \in \mathscr{L}^{+} \backslash\left\{[P]_{+}\right.$, $\left.[Q]_{+}\right\}$. Therefore also $I \cap c$ and $m \cap c$ are fixed by $\Gamma_{P, Q, c}$. By Lemma 11, $\Gamma_{P . Q . c}$ has a subplane $\mathscr{H}_{2}$ as a set of fixed points. Let $n_{2}$ be the order of $\mathscr{M}_{2}$ and let $c^{*}$ be the set of points left fixed by $\Gamma_{P, Q, c}$. With $\mathscr{B}=\left\{c^{* y} \mid \gamma \in \Gamma_{c}\right\}$ we get a $2-\left(n+1, n_{2}+1,1\right)$ design on $c$ (see [13]). The number of blocks through a point is $n / n_{2}=3^{b} / n_{2}$. Hence $n_{2}=3^{d}$ for some $d \in \mathbb{N}$. The total number of blocks equals $(n+1) n /\left(n_{2}+1\right) n=\left(3^{b}+1 / 3^{d}+1\right) \cdot 3^{b-d}$. Hence $b / d \in 2 \mathbb{N}+1$. Since $b$ is even, $d$ is even so $10 \leqslant n_{2}+1=3^{d}+1 \equiv 2(\bmod 4)$. However, $\Gamma_{c^{*}}=N_{\Gamma_{c}}\left(\Gamma_{P, Q, c}\right)$ is sharply 2-transitive on the $n_{2}+1$ points of $c^{*}$, and so $n_{2}+1$ is a power of 2 . This was our final contradiction.

## REFERENCES

1. Bender, H.: 'Endlich zweifach transitive permutationsgruppen deren Involutionen keine Fixpunkte haben', Math. Zeitschr. 104 (1968) 175-204.
2. Benz, W.: Vorlesungen über Geometric der Algebren. Springer-Verlag, Berlin, New York, 1973
3. Bruen, A. and Levinger, B.: 'A Theorem on Permutations of Finite Field', Can. J. Math. 25 (1973), 1060-1065.
4. Chen, Y. and Kaerlein, G.: 'Eine Bemerkung über endiche Laguerre- und MinkowskiEbenen', Geom. Dedicata 2 (1973), (193-194).
5. Dembowski, P. : Finite Geometries, Springer-Verlag, Berlin, Heidelberg, New York, (1968).
6. Feit, W.: 'On a Class of Doubly Transitive Permutation Groups', III. J. Math. 4 (1960), 170-186.
7. Heise, W. : 'Minkowski-Ebenen gerader Ordnung', J. Geom. 5 (1974), 83.
8. Heise, W. and Karzel, H. : 'Symmetrische Minkowski-Ebenen', J. Geom. 3 (1073), 5-20.
9. Hering, C. : 'Endliche zweifach transitive Möbiusebenen ungerader ordnung', Arch. Math. 18 (1967), 212-216.
10. Huppert, B. : 'Scharf dreifach transitive Permutatuinsgruppen', Arch. Math. 13 (1962), 61-72.
11. Huppert, B. : Endliche Gruppen 1, Springer-Verlag, Berlin, Heidelberg, New York, 1967.
12. Ito, N.: 'On a Class of Doubly Transitive Permutation Groups', Ill. J. Math., 6 (1972), 341-352.
13. Kantor, W. : 2-transitive Designs, Mathematical Centre Tracts 57, Amsterdam, 1974, pp. 44-97.

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14. Meschiari, M. and Quattrocchi, P. : 'Una classificazioni delie strutture di incidenza associate a insieme de sostituzioni strettamente 3-transitive finiti', Atti Sem. mat. fis. Univ. Modena, 24 (1975).
15. Ostrom, T. G. and Wagner, A.: 'On Projective and Affine Planes with Transitive Collineation Groups, Math. Zeitschr. 71 (1959), 186-199.
16. Wagner, A.: 'On Finite Affine Line Transitive Planes', Math. Zeitschr. 87 (1965), 1-11.
17. Wielandt, H. : Finite Permutation Groups, Academic Press, New York, 1964.

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A characterization of the classical unitals
by
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ABSTRACT

A characterization of the classical unitals is given in terms of certain geometrical properties.

KEY WORDS \& PHRASES: Unital, inversive plane, generalized quadrangle, Minkowski plane

1. INTRODUCTION

A unital or unitary block design is a $2-\left(q^{3}+1, q+1,1\right)$ design, i.e. an incidence structure of $q^{3}+1$ points, $q^{2}\left(q^{2}-q+1\right)$ lines, such that each line contains $q+1$ points and any two distinct points are on a unique line. If q is a prime power, the absolute points and non-absolute lines of a unitary polarity of $\mathrm{PG}\left(2, q^{2}\right)$ form a unital (see [2]). These unitals are called classical.

In [6], $0^{\prime}$ NAN showed that a classical unital satisfies the following condition.
(I) No four distinct lines intersect in six distinct points (see Figure 1).


Fig. 1

In [5], PIPER conjectured that this property characterizes the classical unitals. Here we shall give a characterization for even $q$ under the assumption that also the following condition holds.
(II) If $L$ is a line, $x$ a point not on $L$, $M$ a line through $x$ meeting $L$ and and $y \neq x$ a point on $M$, then there exists a line $L^{\prime} \neq M$ through $y$ intersecting all lines through $x$ which meet $L$.

To achieve this result we shall give another characterization for all $q$ under the additional assumption that a third condition holds. To formulate this condition we need some notation. If $x$ and $y$ are distinct points, then we denote by $x y$ the line through $x$ and $y$. Given a point $x$, two lines $I$ and $L^{\prime}$ missing $x$ are called $x$-parallel (notation $L^{l l} x^{\prime}$ ) if and only if they intersect the same lines through $x$. Clearly, $\|_{x}$ is an equivalence relation on the set of lines missing $x$, and by (I) and (II), each equivalence class consists of $q$ disjoint lines. Our third condition now reads as follows.
(III) Given a point $x_{r}$, three distinct lines $M_{1}, M_{2}, M_{3}$ through $x$ and points $Y_{i}, z_{i}$ on $M_{i}(i=1,2,3)$ such that $\left(y_{1} y_{2}\right) \|{ }_{x}\left(z_{1} z_{2}\right)$ and $\left(y_{1} y_{3}\right) \|{ }_{x}\left(z_{1} z_{3}\right)$, then also $\left(y_{2} Y_{3}\right) \|_{x}\left(z_{2} z_{3}\right)$.

Clearly, the presence of unitary transvections in PTU (3,q) implies that the classical unitals satisfy conditions (II) and (III). In Section 2 we shall study unitals satisfying (I) and (II). Section 3 is devoted to the proof that unitals satisfying all three conditions are classical. Finally, in Section 4, we shall show that for even $q$, (III) is a consequence of (I) and (II).
2. UNITALS SATISFYING (I) AND (II)

Throughout this section $U$ is a unital on $q^{3}+1$ points with point set $X$ and line set $L$ satisfying ( $I$ ) and (II) above. If $x \in X$, then we denote by $L^{x}$ the set of lines incident with $x$, and $L_{x}$ will be the set of lines missing $x$. Furthermore, $\mathcal{C}_{x}$ will stand for the set of $\|_{x}$-equivalence classes on $L_{x}$. From [1] it is clear that we want to show that the incidence structure which has $L^{X}$ as the set of points, $C_{x}$ as the set of blocks and LIC ( $L \in L^{x}, C \in C_{x}$ ) iff $L$ meets one (hence all) lines of $C$, is the residual of an inversive plane of order $q$. We denote this incidence structure by $I^{*}(x)=\left(L^{x}, C_{x}\right)$. Clearly, $I^{*}(x)$ is a $2-\left(q^{2}, q+1, q\right)$ design.

LEMMA 1. If $\mathrm{x} \in \mathrm{X}$ and $\mathrm{L}, \mathrm{L}^{\prime} \in \mathcal{L}_{\mathrm{x}}$ such that L and $\mathrm{L}^{\prime}$ both meet three distinct lines $M_{1}, M_{2}, M_{3} \in L^{x}$, then $L L_{X} L^{\prime}$, i.e. three distinct points of $I^{*}(x)$ are in at most one block of $I^{*}(x)$.

PROOF. Let $y \in M_{1} \cap L^{\prime}$ and let $L^{\prime \prime}$ be the line through $Y$ which is $x$-parallel to $L$, then $L^{\prime} \neq L^{\prime \prime}$ contradicts (I). $\square$

If $M$ and $M^{\prime}$ are two distinct lines through a point $x$, then an easy counting argument shows that there are $q-2$ lines $N_{1}, \ldots, N_{q-2}$ through $x$ such that no line of $L_{x}$ meets $M, M^{\prime}$ and $N_{i}, i=1, \ldots, q-2$. Put $C^{*}\left(M, M^{\prime}\right):=\left\{M_{1} M^{\prime}\right\} \quad U$ $\left\{N_{1}, N_{2}, \ldots, N_{q-2}\right\}$. We have to show that the $C^{*}(M, M$ ) correspond to circles which will make $I^{*}(x)$ into an inversive plane. We have to show that $N, N^{\prime} \in C^{*}\left(M, M^{\prime}\right) \Rightarrow C^{*}\left(M, M^{\prime}\right)=C^{*}\left(N, N^{\prime}\right)$. Clearly, $C^{*}\left(M, M^{\prime}\right)=C^{*}\left(M^{\prime}, M\right)$ and so it suffices to show that $M^{\prime \prime} \in C^{*}\left(M, M^{\prime}\right) \Rightarrow C^{*}\left(M, M^{\prime}\right)=C^{*}\left(M, M^{\prime \prime}\right)$. This is the
contents of the next lemma.

LEMMA 2. Fix a line $M \in L$ and two distinct points $x$ and $y$ in $M$. For $\left.M^{*},^{\prime \prime} \in L^{X} \backslash M\right\}$ write $M^{\prime} \sim M^{\prime \prime}$ iff no Iine of $L^{Y} \backslash\{M\}$, intersects both $M^{\prime}$ and $M^{\prime \prime}$ or $M^{\prime}=M^{\prime \prime}$. Then $\sim$ is an equivalence relation on $L^{X} \backslash(M)$.

PROOF. For $u, v \in M$ let $A^{*}(u, v)$ be the incidence structure with $L^{u} \backslash\{M)$ as points, $\left.L^{V} \backslash M\right\}$ as lines, and incidence defined by PTB iff $P$ and $B$ meet $\left(P \in L^{u} \backslash\{M\}, B \in L^{V} \backslash\{M\}\right)$. If $u, v$, $w$ are distinct points of $M$, then clearly the mapping $\tau_{v, w}: A^{*}(u, v) \rightarrow A^{*}(u, w)$ defined by

$$
\begin{aligned}
& \widetilde{\tau}_{v, W}^{u}(P):=P, \quad P \in L^{u} \backslash\{M\}, \\
& \sim_{v, w}^{u}(B):=u \text {-parallel of } B \text { through } w, \quad B \in L^{v} \backslash\{M\},
\end{aligned}
$$

is an isomorphism of $A^{*}(u, v)$ onto $A^{*}(u, w)$. Now fix $x, y \in x, x \neq y$. If $q>2$ and $u, v$ are distinct points in $M, u, v \neq x, y$, then

$$
\tilde{\delta}_{u, v}^{x, y}:=\tilde{\tau}_{v, x}^{Y} \tilde{\tau}_{u, y}^{v} \tilde{\tau}_{x, v}^{u} \tilde{\tau}_{y, u^{\prime}}^{x}
$$

is an automorphism of $A^{*}(x, y)$.
Now we claim that
(i) For all $u, v \in M \backslash\{x, y\}, u \neq v$ and for all $\left.P \in L^{x} \backslash M\right\}, \tilde{\delta}_{u, v}^{x}(P) \neq P$ and $\hat{\delta}_{u, v}^{x, y}(P) \sim P$.
(ii) For all $u, v, v^{\prime} \in M \backslash\{x, y\}, u \neq v \neq v^{\prime} \neq u$ and for all $p \in L^{x} \backslash\{M\}$, $\tilde{\delta}_{u, v}^{x, y}(P) \neq \tilde{\delta}_{u, v}^{x}, y(P)$ and $\tilde{\delta}_{u, v}^{x, y}(P) \sim \tilde{\delta}_{u, v}^{x, y}(P)$.

To prove these claims, write $u P v$ for the $u$-parallel of $P$ incident with v. Then $\tilde{\delta}_{u, v}(P)=y(u P v) x$. Suppose $y(u P v) x \neq p$ or $y(u P v) x=P$. Then there is a line $L$ incident with $x$ intersecting $p$ and $y(u p v) x$. Then $L$ intersects uPv in a point a, say. Since au intersects $P$, we now have an o'Nan configuration on the lines $M, P, L$ and au, contradicting (I).

Suppose $y(u P v) x \not y\left(u P v^{\prime}\right) x$ or $y(u P v) x=y\left(u P v^{\prime}\right) x$. Let $L$ be the line through $y$ intersecting $y(u P v) x$ and $y\left(u P v^{\prime}\right) x$. Then $L$ intersects $u P v$ and $u P v^{\prime}$ in points a and $a^{\prime}$, say. Since au'intersects $u P v^{\prime}$, we have an $O^{\prime} N a n$ configuration on $M, L, u P v^{\prime}$ and $a u$, again in contradiction with (I).

For a given $P \in L^{X} \backslash\{M\}$, there are $q-2 \quad Q \in L^{X} \backslash\{M\}, Q \neq P$ such that $Q \sim P$. Fixing $u$ we can make $q-2$ choices for $v \in M \backslash\{x, y, u\}$. Thus, each $Q \in L^{X} \backslash\{M\}$,
 then $Q \sim Q^{\prime}$ by (ii).

Lemma 2 and its proof have a number of important corollaries.

COROLLARY 3. Let $\mathrm{x} \in \mathrm{X}$ and let ${ }_{\mathrm{x}}^{\infty}$ be a new symbol. Put

$$
C^{X}:=\left\{C^{*}\left(M, M^{\prime}\right) \cup\left\{\infty_{X}\right\} \mid M, M^{\prime} \in L^{X}, M \neq M^{\prime}\right\}
$$

Then $I(x):=\left(L^{x} \cup\left\{\infty_{x}\right\}, C^{x} \cup C_{x}\right)$ is an inversive plane of order $q$ with point set $L^{x} \cup\left\{\infty_{x}\right\}$ and block set $C^{x} \cup C_{x}$ and incidence defined in the obvious way.

PROOF. See the discussion preceding Lemma 2. $\quad \square$

COROLLARY 4. For $\mathrm{x}, \mathrm{Y} \in \mathrm{X}, \mathrm{x} \neq \mathrm{y}$, the incidence structure $\mathrm{A}^{*}(\mathrm{x}, \mathrm{y})$ of Lenma 2 is isomorphic to the derived design $I(x)^{x y}$ with ${ }_{x}{ }_{x}$ and the lines through $\infty_{x}$ removed. The affine plane $I(x)^{x y}$ admits a dilatation group of order $q-1$ with centre ${ }^{\infty}{ }_{x}$.

PROOF. The automorphisms $\tilde{\delta}_{u, v}^{x}, y$ of $A^{*}(x, y)$ induce $q-2$ distinct nonidentity dilatations with centre $\infty_{x}$ on $I(x)^{x y}$. Since $I(x)^{x y}$ has order $q$, these are the non-identity elements of the dilatation group with centre ${ }^{\infty} \mathrm{x}$ of order $q-1 . \square$

COROLLARY 5. Let $\mathrm{L} \in \mathrm{L}$ and let $\mathrm{x}_{1}, \mathrm{x}_{2}, \ldots, \mathrm{x}_{\mathrm{q}+1}$ be the points on L. It is possible to partition the set of lines which meet $L$ into classes $A_{i j}$, $1 \leq i, j \leq q+1$, such that for all i and $j$
(i) $\quad\left|A_{i j}\right|=q-1$
(ii) $M \in \mathbf{A}_{\mathbf{i j}} \Rightarrow \mathbf{x}_{\mathbf{i}} \in M$,
(iii) every point $x \in X \backslash L$ is on exactly one line of $U_{k} A_{k j}$,
(iv) no line which meets $L$ in a point $\neq x_{i}$, meets two lines of $A_{i j}$,
(v) for all $k, i^{\prime}, 1 \leq k, i^{\prime} \leq q+1, k \neq i, i^{\prime}$ and for all $M \in A_{i j}$, the $x_{k}-$ parallel of $M$ through $x_{i}$, is in $A_{i}{ }_{j}$ '
(vi) if $1 \leq i^{\prime} \leq q+1, i^{\prime} \neq i$ and $M \in A_{i j}, M^{\prime} \in A_{i}, j$ then there exists $a$ unique $k \in\{1, \ldots, q+1\}$ such that $M$ and $M$ ' are $x_{k}$-parallel.

PROOF. Consider $I\left(x_{1}\right)$. Number the circles of $I\left(x_{1}\right)$ through the two points ${ }^{\infty} x_{1}$ and $L$ of $I\left(x_{1}\right)$ from 1 upto $q+1$. Apart from ${ }^{\infty} x_{1}$ and $L$ each such circle contains $(q-1)$ lines through $x_{1}$. These will be the sets $A_{1, j}, j=1, \ldots, q+1$. For $i>1$ and $1 \leq j \leq q+1$, let $A_{i, j}$ consist of the $(q-1)$ lines through $x_{i}$ which in $I\left(x_{1}\right)$ correspond to the ( $q-1$ ) circles (not through os ${ }_{1}$ ) in the pencil with carrier $L$ and which contains circle $j$ through ${ }^{\infty} x_{1}$ and $L_{\text {. N }}$ Now (i) and (ii) are trivially satisfied. For (iii), note that the $q+1$ lines $x x_{i}$, $i=1, \ldots, q+1$ are in $A_{i j}$ 's with distinct $j$ since the circles in a pencil with carrier L partition the set of points $\neq L$ of $I\left(x_{1}\right)$. To prove the other cases, observe that our subdivision of the set of lines meeting $L$ into the classes $A_{i j}$ would have remained the same if we had started by considering $I\left(x_{i}\right), i>1$ instead of $I\left(x_{1}\right)$. Thus, to prove (iv), it suffices to show that no line $M \in L^{X_{1}} \backslash\{L\}$ can intersect two distinct lines $N_{1} ; N_{2} \in A_{i j}$ with $i>1$. This follows at once, since $N_{1}$ and $N_{2}$ correspond to tangent circles in $I\left(x_{1}\right)$. Also (v) is clear if we take $k=1$ for then $M$ and the $x_{k}$-parallel of $M$ through $x_{i}$, represent the same circle in $I\left(x_{1}\right)$. Finally (vi) follows from (i) , (v) and the easily shown fact that two $l$ ines which meet $L$ cannot be $x_{k}-$ and $x_{\ell}$-parailel for distinct $k$ and $\ell . \square$

Following PIPER [5], we are now able to associate with each line $L$ of $U$ an incidence structure $G Q(E)$ as follows. The points of $G Q(L)$ are the points $x \in X \backslash L$ and the sets $A_{i j}, 1 \leq i, j \leq q+1$. The Lines of $G 2(L)$ are the lines $M$ of $U$ meeting $I$, and $2(q+1)$ new lines, $A_{1}, A_{2}, \ldots, A_{q+1}, B_{1}, B_{2}, \ldots, B_{q+1}$. Incidence in $G Q(L)$ is defined as displayed in the following table.


Incidence in $02(1)$
theorgm 6. Let $U=(X, L)$ be a unital with $q+1$ points on a line satisfying (I) and (II). Then for each line $L \in L, G Q(L)$ is a generalized quadrangle with $q+1$ points on a line and $q+1$ lines through a point. Moreover, any two nonintersecting lines $m_{1}$ and $m_{2}$ of $G Q(L)$ form a regular pair (in the sense of [7]) provided $m_{1}$ and $m_{2}$ do not correspond to lines $M_{1}$ and $M_{2}$ of $U$ such that $M_{1} \in A_{i j}$ and $M_{2} \in A_{k \ell}$ with $i \neq k$ and $j \neq \ell$. In particular, the lines $A_{1}, \ldots, A_{q+1}, B_{1}, \ldots, B_{q+1}$ of $G Q(L)$ are regular.

PROOF. Stxaightforward verification.

We shall see in Section 4 that if all lines of $G 2(L)$ are regular, then U is classical.
3. UNITALS SATISEYING (I), (II) AND (III)

Let $U=(X, L)$ be a unital satisfying (I), (II) and (III). Using (III) it is easy to see that for any three distinct points $x, y, z$ on a line $L$ there is a unique automorphism $T_{y, z}$ of $U$ fixing $x$ and all lines through $x$ and mapping $y$ onto $z$ : if $u \& L$ then $\tau_{y_{r}}^{x}(u)$ is defined to be the point of intersection of $x u$ and the $x$-parallel of yu through $z$, if $v \in L \backslash\{x\}$, fix a point $u \notin L$ and define $\tau_{y, z}^{X}(v)$ to be the point of intersection of $L$ and the $x$-parallel of uv through ${ }^{T} y, z(u)$.

THEOREM 7. Let $U=(X, L)$ be a unital with $q+1$ points on a line satisfying (I), (II) and (III), and let $G$ be the automorphism group of $U$ generated by the $\tau_{y, z}^{x}$. Then $U$ is classical, $G$ is isomorphic to PSU(3, $\left.q^{2}\right)$ and acts on $U$ in the usual way.

PROOF. Clearly $G$ is transitive on $X$. We claim that $G$ acts 2-transitively on $X$ if $q>2$ (the case $q=2$ is left to the reader). To prove this, note that the mappings $\tilde{\tau}_{\bar{X}, z}$ of Lemma 2 are induced by the automorphisms $\tau_{Y, z}^{X}$ of $U$. Hence, also the mappings $\tilde{\delta}_{u, v}^{X, Y}$ of Lemma 2 are induced by automorphisms $\delta_{u, v}^{x, Y} \in G$ of $U$. Since incidence in the inversive plane $I(x)$ is determined by incidence in $U, \delta_{u, v}^{x, y}$ induces an automorphism of $I(x)$. By corollary 4 , this is a dilatation of $I(x)^{x y}$ with centre ${ }^{x} x^{*}$. Therefore, it can also be viewed as a dilatation of $I(x)^{\infty} x$ with centre $x y$. Thus in the affine plane $I(x){ }^{\infty} x$,
each point is the centre of a dilatation. Hence $I(x){ }^{\infty} x$ is a translation plane and the group generated by the dilatations contains the full translation group of $I(x)^{\infty} x([2, p .187])$. Let $T(x)$ be the normal subgroup of $G_{x}$ consisting of elements which induce (possibly identity) translations of $I(x)^{\infty} x$. Then $T(x)$ acts regularly on the points of $I(x)^{\infty} x$, i.e. on $L^{x}$, and for each line $L \in L^{X}, T(x)$ acts regularly on $L \backslash\{x\}$. Thus $T(x)$ is a normal subgroup of $G X$ acting regularly on $X \backslash\{x\}$, and $G$ is 2-transitive. Applying [4] we get that $G$ has a normal subgroup $M$ such that $M \leq G \leq A u t M$ and $M$ acts on $X$ as one of the following groups in its usual 2-transitive representation: a sharply 2 -transitive group, $\operatorname{PSL}\left(2, q^{3}\right), S z\left(q^{3 / 2}\right), \operatorname{PSU}\left(3, q^{2}\right)$, or a group of Ree type. Since $q^{3}+1=(q+1)\left(q^{2}-q+1\right)$ is not a prime power for $q>2$, the first alternative will not occur. If $H \leq G$ and $x, y, z$ are three distinct points of $X$, then the $\mathrm{H}_{\mathrm{xy}}^{3}$-orbit of $z$ is contained in xy , so has length $\leq \mathrm{q}-1$. This excludes $M=\operatorname{PSL}\left(2, q^{3}\right)$ and $M=S z\left(q^{3 / 2}\right)$. Moreover, this argument shows that if $M=\operatorname{PSU}\left(3, q^{2}\right)$ then $U$ is classical, for $M$ has a unique orbit of length $q-1$ on $x \backslash\{x, y\}$, all other orbits have length $\left(q^{2}-1\right) /(q+1,3)([6, p$. 499]). Now the $\tau_{y, z}$ can be identified with the unitary transvections and it follows that $G \simeq \operatorname{PSU}\left(3, q^{2}\right)$. Thus we are left with the case that $M$ is a group of Ree type. Since $q=3^{2 a+1}$, $G$ contains an involution $\delta$ fixing at least two points $x, y \in X$ (Corollary 4). By [4], Lemma $3.3(v)$ and (ix), $\delta \in M$ and $\delta$ fixes $q+1$ points. Since $\delta$ is a dilatation on $I(x)^{\infty} x$ these must be the $q+1$ points of $x y$ and so $U$ is nothing but the Ree unital associated with M. Now, for $L \in L,\left\langle\delta>x \operatorname{PSL}(2, q) \simeq M_{L} \unlhd G_{L}\right.$ and so $\left\langle\tau_{y, z}^{X} \mid x, y, z \in L\right\rangle \leq \operatorname{Aut}(\operatorname{PSL}(2, q))=$ $\operatorname{PrL}(2, q)$, which shows that at least one, and hence all, $\tau_{y, z}^{x} \in M$, i.e. $G=M$ of order $\left(q^{3}+1\right) q^{3}(q-1)$. Now for a 3-5ylow group $T(x)$ of $G, T(x) / T(x) L$ ( $x$ on L) is the elementary abelian translation group of $I(x){ }^{\infty} x$. Hence, for the derived group $T(x)^{(1)}$ of $T(x)$ we find $\left|T(x)^{(1)}\right| \leq\left|T(x)_{L}\right|=q$, contradicting Lemma 3.3(iii) of [4].

## 4. MORE CHARACTERIZATIONS

Let $U=(X, L)$ be a unital satisfying (I) and (II). Consider the following two conditions.
(III') Given a point $x$ and three distinct lines $M_{1}, M_{2}, M_{3}$ through $x$ and points $y_{i}, z_{i}$ on $M_{i}(i=1,2,3)$ such that $\left(y_{1} y_{2}\right)\left\|_{x}\left(z_{1} z_{2}\right),\left(y_{1} y_{3}\right)\right\|\left(z_{1} z_{3}\right)$
and one of the lines $\left(y_{i} y_{j}\right)$ or $\left(z_{i} z_{j}\right)$ meets all three of $M_{1}, M_{2}$ and $M_{3}$, then $\left(y_{2} y_{3}\right) \|_{x}\left(z_{2} z_{3}\right)$.
(IV) Given a point $x$ and two distinct lines $M_{1}$ and $M_{2}$ through $x$ and points $y_{1}, y_{3}, z_{1}, z_{3}$ on $M_{1}, y_{2}, y_{4}, z_{2}, z_{4}$ on $M_{2}$ such that $\left(y_{1} y_{2}\right) l_{x}\left(z_{1} z_{2}\right)$, $\left(y_{1} y_{4}\right) \|_{x}\left(z_{1} z_{4}\right)$ and $\left(y_{2} y_{3}\right) \|_{x}\left(z_{2} z_{3}\right)$, then also $\left(y_{3} y_{4}\right) \|{ }_{x}\left(z_{3} z_{4}\right)$.

Clearly, (III) implies (III') and (IV). The converse is also true.

LEMMA 8. Let $U=(X, L)$ be a unital satisfying (I), (II), (III') and (IV), then also (III) holds.

PROOF. Let $x, M_{i}, Y_{i}, z_{i}, i=1,2,3$ be as in (III). Suppose that $M_{1}, M_{2}$ and $M_{3}$ determine a circle in $I(x)$ not containing $\infty_{x}$, i.e. suppose there is a line through $y_{1}$ intersecting $M_{2}$ in $u_{2}$ and $M_{3}$ in $u_{3}$, say. Let $v_{2}\left(v_{3}\right)$ be the point of intersection of the $x$-parallel of $y_{1} u_{2}$ through $z_{1}$ and $M_{2}\left(M_{3}\right)$. Using (III') we find that $\left(u_{2} y_{3}\right)\left\|\|_{x}\left(v_{2} z_{3}\right)\right.$ and $\left.\left(u_{3} y_{2}\right)\right\|{ }_{x}\left(v_{3} z_{2}\right)$. Hence by (IV), $\left(y_{2} y_{3}\right) \| x\left(z_{2} z_{3}\right)$ and (III) is shown to hold in this case. The remaining case is where $M_{1}, M_{2}$ and $M_{3}$ are on a circle of $I(x)$ containing ${ }^{\infty} x^{\prime}$, i.e. no line of $L_{x}$ meets all three of $M_{1}, M_{2}$ and $M_{3}$. Since the two circles of $I(x)$ corresponding to $y_{1} y_{2}$ and $y_{1} Y_{3}$ cannot be tangent (for otherwise $y_{1} y_{2}=y_{1} y_{3}$ and there is a line intersecting $M_{1}, M_{2}$ and $M_{3}$ ), there is a line $M_{4}$ through $x$ which meets $y_{1} Y_{2}$ in $Y_{4}$ and $z_{1} z_{2}$ in $z_{4}$, say, and which also meets $y_{1} y_{3}$ and $z_{1} z_{3}$. Now looking at $M_{1}, M_{3}$ and $M_{4}$ are applying (III') we see that $\left(y_{3} y_{4}\right) \|_{x}\left(z_{3} z_{4}\right)$. Since $M_{2}$, $M_{3}$ and $M_{4}$ are not on a circle of $I(x)$ (for otherwise this would be the circle determined by $M_{1}, M_{2}$ and $M_{3}$ ), we can apply the previous case and conclude that $\left(y_{2} y_{3}\right) \|_{x}\left(z_{2} z_{3}\right)$.

The reason for considering (III') and (IV) is that in both cases there is a line $M_{i}$ which is intersected by all lines mentioned in the condition. Thus, both (III') and (IV) have a (no doubt awkward) equivalent formulation ( $\overline{\text { III') }}$ respectively ( $\overline{\text { IV }}$ ) into terms of $G 2\left(M_{i}\right)$. Since the classical unital satisfies (III') and (IV), the classical generalized quadrangle $Q(4, q)$ on the points and lines of a hyperquadric in PG(4,q) must satisfy (III') and (IV). So, conversely, if a unital $U$ satisfying (I) and (II) has the property that $G Q(L)$ is isomorphic to $Q(4, q)$ for each line $L$ of $U$, then $U$ is classical.

THEOREM 9. Let 4 be a unital with q+1 points on a line satisfying (I) and (II). If for each line $L$ of $U, G Q(L) \simeq Q(4, q)$, i.e. if every line of $G Q(L)$ is regular, then $U$ is classical.

We are now in a position to prove that for even $q$, (I) and (II) suffice to characterize $U$.

THEOREM 10. Let $U=(x, 1)$ be a unital with $q+1$ points on a line satisfying (I) and (II). If $q$ is even, then $U$ is classical.

PROOF. Let $L$ be a line of $U$ and let $A_{i j}, 1 \leq i, j \leq q+1$ and $A_{i}, B_{i}, i=1, \ldots$ $\ldots, q+1$ be defined as before. For each $x \in x \backslash I$ put

$$
C(x):=\left\{A_{i j} \mid \exists \text { line } M \in A_{i j} \text { incident with } x\right\}
$$

By Corollary 5, $C(x)$ has exactly one point on each of the lines $A_{i}$ and $B_{i}$. $1=1, \ldots, q+1$. We claim that if $x, y \in X \backslash L, x \neq y$, then $|C(x) \cap C(y)| \leq 2$. First suppose $x y$ is a line meeting $L$, $x y \in A_{i f}$, say, then by Corollary 5 (iv), $C(x) \cap C(y)=\left\{A_{i j}\right\}$. Now consider the case where $x y$ is a line of $U$ not meeting L. Suppose $x_{1}, x_{2}, x_{3}$ are distinct points of $L$ such that $x x_{1}, y x_{1} \in A_{1,1}$,
 circles with the following properties: $y x_{1}, y x_{2}, y x_{3}$ all go through the point $x y$ of $I(x)$ and are tangent to $L$ in respectively $x x_{1}, x x_{2}$ and $x x_{3}$. Since $q$ is even, there is a point $\neq \mathrm{xy}$ of $I(\mathrm{x})$ which is also on the circles $\mathrm{yx} \mathrm{x}_{1}, y x_{2}$, $y x_{3}$, i.e. there is a line $M \neq x y$ through $x$ intersecting $y x_{i}, i=1,2,3$. By Lemma $1, L_{y}{ }_{M}$ and so $x y$ intersects $L$, a contradiction. We have shown that each triple $A_{i_{1}}, j_{1}, A_{i_{2}}, j_{2}, A_{i_{3}}, j_{3}$ with $\left|\left\{i_{1}, i_{2}, i_{3}\right\}\right|=\left|\left\{j_{1}, j_{2}, j_{3}\right\}\right|=3$ is covered at most once by a $C(x)$. Since there are $q^{3}-q C(x)$, each such triple is covered exactly once. Thus, with the $A_{i j}$ as points, the $A_{i}$ and $B_{i}$ as lines and the $C(x)$ as circles, we have obtained a Minkowski plane $M(L)$ of even order q. By [3], $M(L)$ is isomorphic to the geometry of points, lines and plane sections of a quadric of index two in PG(3,q). Since $G Q(L)$ is determined by $M(L)$ (the points of $G Q(L)$ correspond to the points and circles of $M(L)$, the lines of $G Q(L)$ correspond to the lines and pencils of $M(L)$, etc.) $G Q(L)$ is isomorphic to $Q(4, q)$ and so $U$ is classical.

## REFERENCES

[1] BUEKENHOUT, F., Existence of unitals in finite translation planes of order $q^{2}$ with a kernel of order $q$, Geometriae Dedicata $\underline{5}$ (1976) 189-194.
[2] DEMBOWSKI, P., Finite geometries, Springer-Verlag, Berlin etc. (1968).
[3] HEISE, W., Minkowski-Ebenen gerader Ordnung, J. Geometry 5 (1974) 83.
[4] HERING, C., W. KANTOR \& G. SEITZ, Finite groups with a split BN-pair of rank 1. I, J. of Algebra 20 (1972) 435-475.
[5] PIPER, F., Unitary block designs, in: Graph Theory and Combinatorics, R.J. Wilson (ed.), Research Notes in Mathematics 34.
[6] O'NAN, M., Automorphisms of Unitary Block Designs, J. of Algebra 20, 495-511 (1972).
[7] THAS, J.A., Combinatorics of partial geometries and generalized quadrangles, in: Higher Combinatorics, Proc. of the Nato Advanced Study Institute, Berlin, Sept. 1-10, 1976, M. Aigner (ed.).

A characterization of two classes of semi partial geometries by their parameters
by
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ABSTRACT

We show that, under mild restrictions on the parameters, semi-partial geometries with $\mu=\alpha^{2}$ or $\mu=\alpha(\alpha+1)$ are determined by their parameters.

KEY WORDS \& PHRASES: Semi-partial geometry, partial geometry, strongly regulas graph

## 1. INTRODUCTION

Let $X$ be a (finite) nonempty set and $L$ a set of subsets of $X$. Elements of $X$ are called points, elements of $L$ are called lines. The pair ( $X, L$ ) is called a partial linear space if any two distinct points are on at most one line.

Two distinct points $x$ and $y$ are called collinear if there exists $L \in L$ such that $x, y \in L$, noncolZinear otherwise. Two distinct 1 ines $L$ and $M$ are called concurpent if $|\mathrm{L} \cap \mathrm{M}|=1$.

We write $x \sim y(x \neq y)$ to denote that $x$ and $y$ are collinear (noncollinear). Similarly $L \sim M(L \neq M)$ means $|L \cap M|=1(|L \cap M|=0)$.

If $x \sim y(L \sim M)$ we denote by $x y(L M)$ the line (point) incident with $x$ and $y$ ( $L$ and $M$ ).

For a nonincident point-line pair ( $x, L$ ) we define:

$$
\begin{aligned}
& {[L, x]:=\{y \in X \mid y \in L, y \sim x\}} \\
& {[x, L]:=\{M \in L \mid x \in M, E \sim M\}}
\end{aligned}
$$

Given positive integers $s, t, \alpha, \mu$, the partial linear space ( $X, L$ ) is called a semi-partial geometry (s.p.g) with parameters $s, t, \alpha, \mu$ if:
(i) every line contains s+1 points,
(ii) every point is on $t+1$ lines,
(iii) for all $x \in X, L \in L, x \notin L$ we have $|[x, L]| \in\{0, \alpha\}$,
(iv) for all $x, y \in X$ with $x \neq y$ the number of points $z$ such that $x \sim z \sim y$ equals $\mu$.
A semi-partial geometry which satisfies $|[x, L]|=\alpha$ for all $x \in X, L \in L$ with $x \notin L$, or equivalently which satisfies $\mu=\alpha(t+1)$, is also called a partial geometry (p.g).

The point-groph of the partial linear space ( $X, L$ ) is the graph with vertex set $X$, two distinct vertices $x$ and $y$ being adjacent iff $x \sim y$. The point-graph of a semi-partial geometry is easily seen to be strongly regular. Let ( $\mathrm{X}, \mathrm{L}$ ) be a semi-partial geometry.

For $x, y \in X, x \notin y$ we define

$$
[x, y]:=\{L \in L|x \in L,|[L, y]|=\alpha\}
$$

It is easy to see that $\alpha=s+1$ iff any two distinct points are collinear iff ( $X, L$ ) is a Steiner system $S(2, s+1,|X|)$. We shall always assume $s \geq \alpha$, hence noncollinear points exist.

Let $x, y \in X, x \neq y$. Then $\mu=|[x, y]| \alpha$ and $|[x, y]| \geq|[x, L]|=\alpha$ if $L \in[y, x]$. Hence, $\mu \geq \alpha^{2}$ and
(*) $\quad \mu=\alpha^{2} \Leftrightarrow \forall K \in[x, y], L \in[y, x]: K \sim L$, (**) $\quad \mu=\alpha(\alpha+1) \Leftrightarrow$ every line $K \in[x, y]$ intersect every line $L \in[y, x]$ but one.

This is the basic observation we use in showing that, under mild restrictions on the parameters, semi partial geometries with $\mu=\alpha^{2}$ or $\mu=$ $\alpha(\alpha+1)$ satisfy the Diagonal Axiom (D).
(D) : Let $x_{1}, x_{2}, x_{3}, x_{4}$ be four distinct points no three on a line, such that $x_{1} \sim x_{2} \sim x_{3} \sim x_{4} \sim x_{1} \sim x_{3}$.
Then also $\mathrm{x}_{2} \sim \mathrm{x}_{4}$.
From Debroey [1], it then follows that such a semi-partial geometry is known.
2. SEMI-PARTIAL GEOMETRIES WITH $\mu=\alpha^{2}$.

Our first theorem deals with the case $\alpha=1, \mu=1$.
THEOREM 1. Every strongly regular graph with parameters $(\mathrm{n}, \mathrm{k}, \lambda, \mu=1)$ is the $\overline{\text { point-graph }}$ of a s.p.g. with $s=\lambda+1, t=\frac{k}{\lambda+1}-1, \alpha=1, \mu=1$.

PROOF: Let ( $\mathrm{X}, \mathrm{E}$ ) be a strongly regular graph with $\mu=1$, and let $\mathrm{x} \in \mathrm{X}$. Since two nonadjacent points in $\Gamma(x)$ cannot have a common neighbour in $\Gamma(x)$, the induced subgraph on $\Gamma(x)$ in the union of cliques. This induced subgraph has valency $\lambda$, so it is the union of $\frac{k}{\lambda+1}$ cliques of size $\lambda+1$.

Next we deal with the case $\alpha=2, \mu=4$.

TH ZOREM 2. Let (X,L) be a s.p.g. with parameters s,t, $\alpha=2, \mu=4$. Then ( $\mathrm{X}, \mathrm{L}$ ) satisfies ( D ).

PROOF. Let $x_{1}, x_{2}, x_{3}, x_{4}$ be four distinct points no three on a line, such that $x_{1} \sim x_{2} \sim x_{3} \sim x_{4} \sim x_{1} \sim x_{3}$. If $x_{2} \nsim x_{4}$, then we can apply (*) to the points $x_{2}$ and $x_{4} *$ Since $x_{1} x_{4} \in\left[x_{4}, x_{2}\right]$ and $x_{2} x_{3} \in\left[x_{2}, x_{4}\right], x_{1} x_{4}$ and $x_{2} x_{3}$ intersect in a point $\neq x_{2}, x_{3}$. Now $3 \leq\left|\left[x_{1}, x_{2} x_{3}\right]\right| \leq \alpha=2$, a contradiction.

Let $U$ be a set containing $t+3$ elements. Then we denote by $U_{2,3}$ the s.p.g. which has as points the 2 -subsets of U , as lines the 3 -subsets of U together with the natural incidence.

The parameters are $s=2, t, \alpha=2, \mu=4$.
DEBROEY [1] showed that a s.p.g. with $t>1, \alpha=2, \mu=4$ satisfying (D) is isomorphic to a $U_{2,3}$. Hence we have the following theorem.

THEOREM 3. A s.p.g. with $\mathrm{t}>1, a=2, \mu=4$ is isomorphic to $a \mathrm{U}_{2,3}$. $\square$
REMARK. A s.p.g. with $t=1, \alpha=2, \mu=4$ is isomorphic to the geometry of edges and vertices of the complete graph $\mathrm{K}_{\mathrm{s}+2}$

We now consider the case $\alpha>2$. For the remainder of this sec wet ( $\mathrm{X}, \mathrm{L}$ ) be a s.p.g with $\alpha>2$ and $\mu=\alpha^{2}$.

LEMMA 1. Let $\mathrm{x} \in \mathrm{X}, \mathrm{L} \in \mathrm{L}, \mathrm{X} \in \mathrm{L}$ such that $[\mathrm{L}, \mathrm{X}]=\left\{\mathrm{z}_{1}, \ldots, \mathrm{z}_{\alpha}\right\}$. Let M be a line through $z_{1}$ intersecting $\mathrm{x}_{2}$ in a point $\mathrm{u} \neq \mathrm{x}, \mathrm{z}_{2}$. Suppose there exists $\mathrm{y} \in \mathrm{L}, \mathrm{y} \neq \mathrm{z}_{1}, \ldots, \mathrm{z}_{\alpha}$ with $\mathrm{u} \nmid \mathrm{y}$. Then M intersects $\mathrm{xz} \mathrm{z}_{\mathrm{i}}$ for all $\mathrm{i}=1, \ldots, \alpha$ (see figure 1).


Figure 1.

PROOF. By (*) applied to $x$ and $y$, the $\alpha$ lines $L=L_{1}, L_{2}, \ldots, L_{\alpha}$ of $[y, x]$ intersect the $\alpha$ lines $x z_{1}, \ldots, x z_{\alpha}$ of $[x, y]$. In particular $L_{1}, \ldots, L_{\alpha}$ intersect $x z_{2}$. Hence $[y, u]=[y, x]=\left\{L_{1}, \ldots, L_{\alpha}\right\}$.

Since $M \in[u, y], M$ intersects $L_{1}, \ldots, L_{\alpha}$ in points $v_{1}=z_{1}, v_{2}, \ldots, v_{\alpha}$ respectively, If $x \sim v_{i}$ for all $i$, then the $\alpha+1$ points $u, v_{1} v_{2}, \ldots, v_{\alpha}$ on $M$ are all collinear with $x$, a contradiction. Hence $x \not v_{i}$ for some i. Since $L_{i}$ intersects $x z_{1}, \ldots, x z_{\alpha}$ it follows that $\left[x, v_{i}\right]=[x, y]=\left\{x z_{1}, \ldots, x z_{\alpha}\right\}$. Since $M \in\left[v_{i}, x\right]$, $M$ intersects all lines in $\left[x, v_{i}\right]$, $\square$

LEMMA 2. Let $\mathrm{x} \in \mathrm{X}, \mathrm{L} \in \mathrm{L}, \mathrm{x} \notin \mathrm{L}$ such that $[\mathrm{L}, \mathrm{x}]=\left\{\mathrm{z}_{1}, \ldots, \mathrm{z}_{\alpha}\right\}$. Let M be $a$ line through $z_{1}$ intersecting $\mathrm{xz}_{2}$ in a point $\mathrm{u} \neq \mathrm{x}, \mathrm{z}_{2}$. If $\mathrm{s}>\alpha$, then M intersects $\mathrm{xz}_{\mathbf{i}}$, for all $\mathbf{i}=1, \ldots, \alpha$.

PROOF. Assume that $M$ intersects $x_{i}$, $i=1, \ldots, \beta(2 \leq \beta<\alpha)$ in points $u_{1}=z_{1}$, $u_{2}=u, \ldots, u_{\beta}$ respectively and does not intersect $x z_{\beta+1}, \ldots, x z_{\alpha}$. Take $y \in L, y \neq z_{1}, \ldots, z_{\alpha}$. By lemma $1 \mathrm{y} \sim u_{i}, i=1, \ldots, \beta$.

Since $\mid[M, x] \|=\alpha$, there is a $v \in M$ such that $v \sim x, v \neq u_{1}, \ldots, u_{\beta}$. Also $v \sim z$ for all $z \in \underset{i=1}{B}\left[y u_{i}, x\right]$, for if $v \nLeftarrow z$ for some $z \in\left[y u_{i}, x\right]$, then $v x \in[v, z]$ and $y u_{i} \in[z, v]$. Hence $v x \sim y u_{i}$ and so $y u_{i}$ intersects the $\alpha+1$ lines $x v, x z_{1}, \ldots, x z_{\alpha}$ through $x$, a contradiction. The points of $\underset{i=1}{b}\left[y u_{i}, x\right]$ are therefore on the $\alpha$ lines $M=v z_{1}, v z_{2}, \ldots, v z_{\alpha}$ of $[v, y]$.

Since $s>\alpha$ we can take $y^{\prime} \in L$ such that $y^{\prime} \neq y, z_{1}, \ldots, z_{\alpha}$.
Now if $z \in\left[y u_{2}, x\right]$, then $z \sim y^{\prime}$. Indeed, as shown $z$ is on some $v z_{i}$ and since $v z_{i}$ intersects at most $\alpha-1$ of the lines $\mathrm{xz}_{1}, \ldots, \mathrm{xz} z_{\alpha}$, it follows from Lemma 1 that every point of intersection of $v z_{i}$ and a line $x z_{j}$, so in particular $z$, is collinear with $y^{\prime}$.

But now we have $\left|\left[y u_{2}, y^{\prime}\right]\right| \geq\left|\left[\mathrm{yu}_{2}, x\right] \cup\{y\}\right|=\alpha+1$, a contradiction.
LEMMA 3. Let $\mathrm{x} \in \mathrm{X}, \mathrm{L} \in \mathrm{L}, \mathrm{x} \notin \mathrm{L}$ such that $[\mathrm{L}, \mathrm{x}]=\left\{\mathrm{z}_{1}, \ldots, \mathrm{z}_{\alpha}\right\}$. If $\mathrm{s}>\alpha$, then every line $M$ not through $x$ which intersects two lines of $[x, L]=$ $\left\{x z_{1}, \ldots, x z_{\alpha}\right\}$ also intersects L and all tines of $[\mathrm{x}, \mathrm{L}]$.

PROOF. The number of pairs $(u, v) \neq\left(z_{1}, z_{2}\right)$ such that $u \in x z_{1}, v \in x z_{2}$, $u, v \neq x, u \sim v$ equals $s(\alpha-1)-1$. Every line $M \neq x z_{1}, \ldots, x z_{\alpha}$ which intersects L and $x z_{1}, \ldots, x z_{\alpha}$ gives rise to such a pair ( $u, v$ ). By (*) and lemma 2 the number of these lines equals $(s+1-\alpha)(\alpha-1)+\alpha(\alpha-2)=s(\alpha-1)-1$. $\square$

Let $L_{1}, L_{2} \in L$ intersect in a point $x$. If $L$ is any line intersecting $L_{1}$ and $\mathrm{L}_{2}$ not in x , we let $\mathrm{L}_{3}, \mathrm{~L}_{4}, \ldots, \mathrm{~L}_{\alpha}$ be the other lines in $[\mathrm{x}, \mathrm{L}]$. By lema 3, $\mathrm{L}_{3}, \mathrm{~L}_{4}, \ldots, \mathrm{~L}_{\alpha}$ are independent of the choice of L . Put

$$
\begin{aligned}
& L\left(L_{1}, L_{2}\right):=\left\{L_{1}, L_{2}, \ldots, L_{\alpha}\right\} \cup\left\{L \in L \mid L \sim L_{1}, L_{2}, L L_{1} \neq x \neq L L_{2}\right\}, \\
& x\left(L_{1}, L_{2}\right):=\underset{L \in L\left(L_{1}, L_{2}\right)}{U} \quad
\end{aligned}
$$

LEMMA 4. Let $\mathrm{L}_{1}, \mathrm{~L}_{2} \in L, \mathrm{~L}_{1} \sim \mathrm{~L}_{2}$. If $\mathrm{s}>\alpha$, then $\left\langle\mathrm{L}_{1}, \mathrm{~L}_{2}\right\rangle:=\left(\mathrm{X}\left(\mathrm{L}_{1}, \mathrm{~L}_{2}\right), L\left(\mathrm{~L}_{1}, \mathrm{~L}_{2}\right)\right)$ is a partial geometry (in fact a dual design) with parameters $\widetilde{s}=s$, $\tilde{t}=\alpha-1, \tilde{\alpha}=\alpha$.

PROOF. Clearly two points are on at most one line and each line contains s+1 points. Using (*) and Lemma 3 it follows immediately that every point $x \in X\left(L_{1}, L_{2}\right)$ is on $\alpha$ lines of $L\left(L_{1}, L_{2}\right)$ so $\tilde{t}+1=\alpha$. It also follows immediately that any two lines of $L\left(L_{1}, L_{2}\right)$ intersect, hence $\tilde{\alpha}=\tilde{t}+1=\alpha$.

Notice that for $M_{1}, M_{2} \in L\left(L_{1}, L_{2}\right), M_{1} \neq M_{2}, M_{1} \sim M_{2}$ we have $\left\langle M_{1}, M_{2}\right\rangle=$ $\left\langle L_{1}, L_{2}\right\rangle$. Notice also that for any two noncollinear points $x$ and $y$ of $\left\langle L_{1}, L_{2}\right\rangle$ there are $\tilde{\mu}=\tilde{\alpha}(\tilde{t}+1)=\alpha^{2}=\mu$ points $z \in X\left(L_{1}, L_{2}\right)$ collinear with both $x$ and $y$, i.e. the common neighbours of $x$ and $y$ in ( $x, L$ ) are the common neighbours of $x$ and $y$ in $<L_{1}, L_{2}>$.

THEOREM 4. Let ( $\mathrm{X}, \mathrm{L}$ ) be a s.p.g. with parameters $s, t, \alpha(>2), \mu=\alpha^{2}$. If $s>\alpha$ and $t \geq \alpha$, then ( $\mathrm{X}, \mathrm{L}$ ) satisfies ( D ).

PROOF. Let $x_{1}, x_{2}, x_{3}, x_{4}$ be four distinct points no three on a line, such that $x_{1} \sim x_{2} \sim x_{3} \sim x_{4} \sim x_{1} \sim x_{3}$.

Suppose $x_{2} \nsim x_{4}$. Since $x_{2} \sim x_{1} \sim x_{4}$ it follows that

$$
x_{1} \in\left\langle x_{4} x_{3}, x_{2} x_{3}\right\rangle
$$

In ( $\mathrm{X}, \mathrm{L}$ ) there are $\lambda=\mathrm{s}-1+(\alpha-1)$ t points collinear with both $x_{1}$ and $x_{3}$. In $\left\langle x_{4} x_{3}, x_{2} x_{3}\right\rangle$ there are $\tilde{\lambda}=\tilde{s}-1+(\tilde{\alpha}-1) \tilde{t}=(s-1)+(\alpha-1)^{2}$ points collinear with both $x_{1}$ and $x_{3}$. Since $t \geq \alpha=\tilde{t}+1$ it follows that $\tilde{\lambda}<\lambda$ and so there exists $x_{5} \in X \backslash X\left(x_{4} x_{3} ; x_{2} x_{3}\right)$ such that $x_{1} \sim x_{5} \sim x_{3}$. Now application of
( +

$$
\begin{array}{ll}
\text { to } x_{1}, x_{5}, x_{3}, x_{4} \text { yields } & x_{5} \sim x_{4}, \\
\text { to } x_{1}, x_{2}, x_{3}, x_{5} \text { yields } & x_{5} \sim x_{2}, \\
\text { to } x_{4}, x_{1}, x_{2}, x_{5} \text { yields } & x_{2} \sim x_{4} .
\end{array}
$$

DEBROEY [1] showed that a s.p.g. with parameters $s, t, \alpha(>2), \mu=\alpha^{2}$ satisfying ( $D$ ) is of the following type: the "points" are the lines of $P G(d, q)$, the "lines" are the planes in $\operatorname{PG}(\mathrm{d}, \mathrm{q})$ for some prime power q and $\mathrm{d} \in \mathbb{N}$, $d \geq 4$. In this case $s=q(q+1), t=(q-1)^{-1}\left(q^{d-1}-1\right)-1, \alpha=q+1, \mu=(q+1)^{2}$. THEOREM 5. Let ( $\mathrm{X}, \mathrm{L}$ ) be a s.p.g. with parameters $\mathrm{s}, \mathrm{t}, \alpha(>2), \mu=\alpha^{2}$. If $\mathrm{s}>\alpha$ and $\mathrm{t} \geq \alpha$, then ( $\mathrm{X}, L$ ) is isomorphio to the s.p.g. consisting of the lines and planes in $\operatorname{PG}(\mathrm{d}, \mathrm{q})$. In particular $\mathrm{s}=\mathrm{q}(\mathrm{q}+1), \mathrm{t}=(\mathrm{q}-1)^{-1}\left(\mathrm{q}^{\mathrm{d}-1}-1\right)-1$, $\alpha=q+1, \mu=(q+1)^{2}$.

The only interesting case remaining is $s=\alpha$. Now if ( $X, E$ ) is a Moore graph of valency $r$, i.e. a strongly regular graph with $\lambda=0, \mu=1$, then ( $\mathrm{X}, \operatorname{\{ r}(\mathrm{x}) \mid \mathrm{x} \in \mathrm{X}\}$ ) is easily seen to be a s.p.g. with parameters $\mathrm{s}=\mathrm{t}=\alpha=$ $=r-1, \mu=(r-1)^{2}$ (here $\Gamma(x)=\{y \in X \mid(x, y) \in E\}$ ). The point graph of this s.p.g. is the complement of (X,E). Such a s.p.g. does not satisfy (D) for $r>2$. From the following theorem follows immediately that a s.p.g. with $\mu=\alpha^{2}, s=\alpha$ is necessarily of this type.

THEOREM 6. Let ( $\mathrm{X}, L$ ) be a s.p.g. with $\mathrm{t} \geq \alpha, \mu=\alpha^{2}$ and $\mathrm{s}=\alpha$. Then $\mathrm{t}=\alpha$.
PROOF, Let $x, y \in X, x \neq y$. Let $[x, y]=\left[L_{1}, \ldots, L_{\alpha}\right\},[y, u]=\left[M_{1}, \ldots, M_{\alpha}\right]$ and put $z_{i j}=L_{i} M_{j}, i, j=1, \ldots, \alpha$ (see figure 2).


Figure 2.

The number of ( $z_{i j}, z_{k \ell}$ ) with $i \neq k, j \neq \ell, z_{i j} \sim_{k \ell}$ equals $\alpha^{2} \cdot(\alpha-1)(\alpha-2)$. Now let $K$ be a line through $x, K \neq L_{1}, \ldots, L_{\alpha}$, and let $u$ be a point on $K$, $\mathrm{u} \neq \mathrm{x}$.

Then $u$ is collinear with ( $\alpha-1$ ) of the $\alpha$ points $z_{i, 1}, \ldots, z_{i, \alpha}$, for $i=1, \ldots, \alpha$. Since $u \neq y, u$ is collinear with all of $z_{1, j}, \ldots, z_{\alpha, j}$ or with none, for $j=1, \ldots, \alpha$.

It follows that there are $\alpha$ lines through $u$ intersecting ( $\alpha-1$ ) of the $\alpha$ lines $M_{1}, \ldots, M_{\alpha}$. Hence each point $u \neq x$ on $K$ gives rise to $\alpha(\alpha-1)(\alpha-2)$ pairs ( $z_{i j}, z_{k \ell}$ ) as described, so $K$ gives rise to all $\alpha^{2}(\alpha-1)(\alpha-2)$ pairs ( $z_{i j}, z_{k \ell}$ ).

Suppose $t>\alpha$, then we can find two such lines $K$ and $K^{\prime}$. It follows that for $u \in K$, the $\alpha$ lines through $u$ intersecting ( $\alpha-1$ ) of the $\alpha$ lines $M_{1}, \ldots, M_{\alpha}$ also intersect $K^{\prime}$. But now $\left|\left[u, K^{\prime}\right]\right|=\alpha+1$, a contradiction.

## 3. SEMI-PARTIAL GEOMETRIES WITH $\mu=\alpha(\alpha+1)$.

In this section ( $X, L$ ) is a semi-partial geometry with parameters $s, t, \alpha$ and $\mu=\alpha(\alpha+1)$.

If $\mathrm{x}, \mathrm{y} \in \mathrm{X}, \mathrm{x} \nmid \mathrm{y}$ we shall always denote the $\alpha+1$ lines in $[\mathrm{x}, \mathrm{y}]$ by $\mathrm{K}_{1}, \ldots, \mathrm{~K}_{\alpha+1}$, and the ( $\alpha+1$ ) lines in $[y, x]$ by $L_{1}, \ldots, \mathrm{~L}_{\alpha+1}$. By ( $* *$ ) we can number these lines in such a way that $K_{i} \cap L_{i}=\emptyset, i=1, \ldots, \alpha+1$ and $K_{i} \cap L_{j} \neq \emptyset, i, j=1, \ldots, \alpha+1, i \neq j$ (see figure 3 ).


Figure 3.
Again our aim will be to show that the diagonal axiom (D) holds. We first
deal with the case $\alpha=2$.

LEMMA 5. If $\alpha=2$ and $t>s$, then a set of 3 collinear points not on one line can be extended to a set of 4 collinear points no 3 on a line.

PROOF. Let $x$, $a$ and $b$ be three distinct collinear points not on one line. There are $t-1$ lines $\neq \mathrm{xa}, \mathrm{ab}$ through a and on each of those lines there is a point $y_{i} \sim b, y_{i} \neq a, i=1, \ldots, t-1$. Suppose $y_{i} \neq x$ for $a 11 i=1, \ldots, t-1$. Now for each $i=1, \ldots, t-1$, $a y_{i} \neq x b$ (for otherwise $|[a, x b]| \geq 3$ ) and $b_{i} \notin x a$. Also xa, xb $\in\left[x, y_{i}\right]$ and $a y_{i}, b y_{i} \in\left[y_{i}, x\right]$. Hence, by (**) there is a third line through $y_{i}$ intersecting $x a$ and $x b$ in points $u_{i}$ and $v_{i}$ respectively. Clearly $u_{i} \neq u_{j}$ if $i \neq j$, for $u_{i}=u_{j}$ implies $x, v_{i}, v_{j} \in\left[u_{i}, x b\right]$. Thus xa contains $t+1>s+1$ points (namely $x ; a_{1} u_{1}, \ldots, u_{t-1}$ ), a contradiction. $\square$

LEMMA 6. Suppose $\alpha=2$. If $\mathrm{x}_{1}, \mathrm{x}_{2}, \mathrm{x}_{3}, \mathrm{x}_{4}$ are fow distinct collineax points, no three on a line, then no point can be collinear with exactly three of these four points.

PROOF. Suppose $x_{5}$ is collinear with $x_{2}, x_{3}, x_{4}$ and $x_{1} \psi x_{5}$. Clearly $x_{5} \& x_{2} x_{3}$, $x_{2} x_{4}, x_{3} x_{4}$. Hence $\left\{x_{1} x_{2}, x_{1} x_{3}, x_{1} x_{4}\right\}=\left[x_{1}, x_{5}\right]$ and $\left\{x_{5} x_{2}, x_{5} x_{3}, x_{5} x_{4}\right\}=\left[x_{5}, x_{1}\right]$ so $x_{5} x_{2}$ has to intersect $x_{1} x_{3}$ or $x_{1} x_{4}$ by ( $* *$ ). But then $\left|\left[x_{2}, x_{1} x_{3}\right]\right|$ or $\left|\left[x_{2}, x_{1} x_{4}\right]\right|>2$, a contradiction. $\square$

LEMMA 7. Same hypothesis as in lema 6. Then the only points coltinear with exactly two points of $\left\{x_{1}, x_{2}, x_{3}, x_{4}\right\}$ are the points on the lines $x_{1} x_{j}, i \neq j$. PROOF. Suppose $x_{5} \sim x_{1}, x_{4}$ and $x_{5} \not \subset x_{2}, x_{3}, x_{5} \notin x_{1} x_{4}$ (see figure 4).


Figure 4.

Apply ( $* *$ ) to $x_{3}$ and $x_{5}$ to get a line ab through $x_{3}$ with $a \in x_{5} x_{4}, b \in x_{5} x_{1}$. Similarly (**) applied to $x_{5}$ and $x_{2}$ gives us a line cd through $x_{2}$ with $c \in x_{5} x_{4}, d \in x_{5} x_{1}$. Clearly $b \not \& c$ so we can apply (**) to $b$ and c. It follows that $a b \cap c d=\emptyset$. Also $x_{2} \nmid a$ and ( $* *$ ) applied to $x_{2}$ and a yields: $a b \cap c d \neq \emptyset$ or $a b \cap x_{2} x_{4} \neq \emptyset$. Hence $a b \cap x_{2} x_{4} \neq \emptyset$, a contradiction since $\left\{x_{2}, x_{4}\right\}=\left[x_{2} x_{4}, x_{3}\right]$.

THEOREM 7. If $(\mathrm{X}, \mathrm{L})$ is a s.p.g with parameters $\mathrm{s}, \mathrm{t}, \alpha=2, \mu=6$ and $\mathrm{t}>\mathrm{s}$, then ( $\mathrm{X}, \mathrm{L}$ ) satisfies ( D ).

PROOF. Let $x_{1}, x_{2}, x_{3}$ and $x_{4}$ be four distinct points no three on a line such that $x_{4} \sim x_{1} \sim x_{2} \sim x_{3} \sim x_{4} \sim x_{2}$. By Lemma 5 there exists $x_{5} \sim x_{2}, x_{3}, x_{4}$.

By Lermas 6 and $7 x_{1} \sim x_{3}, x_{5}$.
REMARK. If ( $\mathrm{X}, \mathrm{L}$ ) is a s.p.g but not a partial geometry, then $\mathrm{t} \geq \mathrm{s}$ (see DEBROEY \& THAS [2]). Using the integrality conditions for the multiplicities of the eigenvalues of a strongly regular graph it follows that a s.p.g with $s=t, \alpha=2$ and $\mu=6$ satisfies $\left(8 s^{2}-24 s+25\right) \mid\left\{8(s+1)\left(2 s^{3}-9 s^{2}+19 s-30\right)\right\}^{2}$. From this one easily deduces an upper bound for $s$. The remaining cases were checked by computer and only $s=t=28$ survived. Thus, every s.p.g which is not a partial geometry satisfies (D) or has $s=t=28$ (and 103125 points).

We now turn to the case $\alpha \geq 3$. We shall make two additional assumptions in this case. The first assumption is $\alpha \neq 3$, the second assumption is $s \geq f(\alpha)$ where $f$ is defined in Lemma 9. Notice that this bound on $s$ is used only in the proof of Lemma 9.

LEMMA 8. Let $\mathrm{x}, \mathrm{y} \in \mathrm{x}, \mathrm{x} \not \mathrm{y}$ and suppose $[\mathrm{x}, \mathrm{y}]=\left[\mathrm{K}_{1}, \ldots, \mathrm{~K}_{\alpha+1}\right],[\mathrm{y}, \mathrm{x}]=$ $\left[L_{1}, \ldots, L_{\alpha+1}\right]$ such that $K_{i} \cap L_{i}=\emptyset, i=1, \ldots, \alpha+1$. If $M$ is a line inter. secting $\sigma \geq 1$ lines of $[x, y], \tau \geq 1$ lines of $[y, x]$ and $\sigma<\tau$, then $\sigma=\alpha-1$ and $\tau=\alpha$.

PROOF. Since $\sigma<r$, there exists a point of intersection $u$ of $M$ with a line $L_{i} \in[y, x]$ such that $u$ is not on one of the lines of $[x, y]$. Then $u \neq x$ and so, applying ( $* *$ ) to $u$ and $x$, it follows that $M \in[u, x]$ intersects $\alpha-1$ of the $\alpha$ lines $K_{1}, K_{2}, \ldots, K_{i-1}, K_{i+1}, \ldots, K_{\alpha+1} \in[x, u]$. Thus $\alpha-1 \leq \sigma<\tau \leq \alpha$, which proves our claim. $\square$

LEMMA 9. Let $\mathrm{X} \in \mathrm{X}$ and $\mathrm{L} \in \mathrm{L}$ such that $\mathrm{x} \notin \mathrm{L}$ and x is collinear with a points $z_{2}, z_{3}, \ldots, z_{\alpha+1}$ on L . Let M be a line through $\mathrm{z}_{\alpha+1}$ meeting $\mathrm{x} z_{\alpha}$ in a point $u \neq x, z_{\alpha}$. Suppose $\mathrm{s} \geq f(\alpha)$ where $f(4)=12, f(5)=16, f(6)=f(7)=17$, $f(8)=18, f(9)=19, f(10)=21, f(11)=23, f(\alpha)=2 \alpha(\alpha \geq 12)$. Then $M$ inter sects at least $\alpha-1$ lines of $[\mathrm{x}, \mathrm{L}]$.

PROOF. Suppose $M$ does not meet at least two lines of $[x, L], x z_{2}$ and $x z_{3}$, say. Since $s \geq 2 \alpha$ we can find $y \in L$ such that $x \not f y \%$. Let $[x, y]=\left\{K_{1}, K_{2}=x z_{2}, \ldots\right.$, $\left.K_{\alpha+1}=x Z_{\alpha+1}\right\}$ and $[y, x]=\left\{L_{1}=L, L_{2}, L_{3}, \ldots, L_{\alpha+1}\right\}$ with $K_{i} \cap L_{i}=\emptyset$.

Looking at $u$ and $y$ we find that $M$ intersects $\alpha-1$ of the $\alpha$ lines $L_{i}$, $i \neq \alpha$. Every point $L_{i} M$ which is collinear with $x$ is on a line $K_{j}, j \neq \alpha$. If $L_{i} M \sim x$ for these $\alpha-1 i ' s$, we find that $M$ meets at least $\alpha$ of the lines $K_{1}, \ldots, K_{\alpha+1}$, hence at least $\alpha-1$ of the lines $K_{2}, \ldots, K_{\alpha+1}$, a contradiction. Let $t=L_{i} M$ be a point not collinear with $x$. Considering $x / t$ we see that $M$ intersects $\alpha-1$ of the $\alpha$ lines in $[x, y] \backslash\left\{K_{i}\right\}$. This shows that $i=2$ or 3 , so there are at most two such points $t$, and that $M$ meets $K_{1}, K_{4}, K_{5}, \ldots, K_{\alpha+1}$. Let $V=\left\{K_{4} M, K_{5} M, \ldots, K_{\alpha} M\right)$ and count pairs ( $\left.y, v\right), y \in L, y \not \subset x, v \in V$, viv. The number of such pairs is at least $(s-\alpha+1)(\alpha-5)$ (first choose $y, s-\alpha+1$ possibilities, then given $y$ we can find $\alpha-3$ points $L_{i} M-x$ as above, possibly one on $K_{1}(y)$, and one is $z_{\alpha+1}$ ), and at most $(\alpha-3)(\alpha-2)$ (first choose $v$, then $y$ ). It follows that for $\alpha>5, s \leq 2 \alpha-1+\left\lfloor\frac{6}{\alpha-5}\right\rfloor$. Let $W=V \cup\left\{q, q^{\prime}\right\}=$ $=\left\{w \in M \mid w^{\sim} x\right\}$ and count pairs $(y, w), y \in L, y \notin x, w \in W, w \sim y$. This yields $(s-\alpha+1)(\alpha-4) \leq(\alpha-3)(\alpha-2)+2(\alpha-1)$, hence $s \leq 2 \alpha+\left\lfloor\frac{8}{\alpha-4}\right\rfloor$ if $\alpha>4$. Above we saw that for any $y \in L$ with xfyfu, $K_{1}=K_{1}(y)$ meets $M$. But if $s+1>\alpha+$ $+(\alpha-2)+2(\alpha-1)=4 \alpha-4$, we can find $y \in L$ such that $y \not \subset x, u, q$ and $q^{*}$, a contradiction. Therefore we have $s<4 \alpha-4$. We now have obtained a contradiction for all $\alpha \geq 4$ and the lemma is proved.

LEMMA 10. Some hypotheses as in Lemma 9. Then M intersects exactly $\alpha-1$ lines of $[x, L]$.

PROOF. Take $y \in L$, $y \not f x$ and let $K_{i}$ and $L_{i}$ be defined as before. Put $K:=K_{\alpha+1}$ and let $A(x, L)$ be the set of lines $\neq K, L$ through $z_{\alpha+1}$ intersecting at least $\alpha-1$ lines of $[x, L], A(y, K)$ the set of lines $\neq K, L$ through $z_{\alpha+1}$ intersecting at least $\alpha-1$ lines of $[y, K]$. Suppose a lines of $A(x, L)$ intersects $\alpha-1$ lines of $[x, L]$ and $b$ lines of $A(x, L)$ intersect $a$ lines of $[x, L]$. Counting
the points $u \sim z_{\alpha+1}$ on $K_{2}, K_{3}, \ldots, K_{\alpha}$, such that $u \neq x, z_{2}, \ldots, z_{\alpha}$ yields $a(\alpha-2)+b(\alpha-1)=(\alpha-1)(\alpha-2)$. Hence $a=0$ and $b=\alpha-2$ or $a=\alpha-1$ and $b=0$. Thus $|\mathrm{A}(\mathrm{x}, \mathrm{L})|=\alpha-2$ or $\alpha-1$ according as every line in $\mathrm{A}(\mathrm{x}, \mathrm{L})$ intersects all lines or all but one line in [ $x, L]$. A similar result holds for $A(y, K)$. Now $A(x, L)=A(y, K)$, for suppose $N \in A(x, L)$ then by Lenma 8 , $N$ intersects at least $\alpha-1$ lines of $[y, x]$, so at least $\alpha-2 \geq 2$ lines of $[y, k]$. Hence $N \in \mathbb{A}(y, K)$ by Lemma 9. Similarly, $N \in A(y, k)$ implies $N \in A(x, L)$. Suppose $|A(x, L)|=\alpha-2$, i.e. there are $\alpha-2$ lines through $z_{\alpha+1}$ intersecting all lines of $[x, L] \cup[y, K]$. It follows that $K_{2} L_{\alpha+1} \notin z_{\alpha+1}$ so we can apply (**) to $K_{2} L_{\alpha+1}$ and $z_{\alpha+1}$. This shows that $L_{\alpha+1} \in\left[K_{2} L_{\alpha+1}, z_{\alpha+1}\right]$ intersects all $N \in A(y, K) \subseteq\left[z_{\alpha+1}, K_{2} L_{\alpha+1}\right]$, a contradiction, for $L_{\alpha+1} \sim N$ implies $|[y, N]| \geq$ $\geq \alpha+1$.

LEMMA 11. Let $\mathrm{x} \in \mathrm{X}, \mathrm{L} \in \mathrm{L}$ such that x is collinear with a points $z_{2}, \ldots, z_{\alpha+1}$ on L. Let $M$ be a line through $z_{\alpha+1}$ intersecting a -1 lines of $[\mathrm{x}, \mathrm{L}]$ and let $\mathrm{y} \in \mathrm{L}$, $\mathrm{y} \nmid \mathrm{x}$. Then, if $[\mathrm{x}, \mathrm{y}]=\left(\mathrm{K}_{1}(\mathrm{y}), \mathrm{K}_{2}=\mathrm{xz} z_{2}, \ldots, \mathrm{~K}_{\alpha+1}=\mathrm{xz}{ }_{\alpha+1}\right\}$, M intersects $\mathrm{K}_{1}(\mathrm{y})$.

PROOF. Suppose $M$ does not intersect $K_{2}$, say. As shown in Lenma $10, M$ also intersects $\alpha-1$ lines of $\left[y, K_{\alpha+1}\right]=\left\{L_{1}=\mathrm{L}, \mathrm{L}_{2}, \ldots, \mathrm{~L}_{\alpha}\right\}$. So $M$ intersects at least one of $\mathrm{L}_{\alpha-1}$ and $\mathrm{L}_{\alpha}$ and since $\alpha \geq 4, \mathrm{~L}_{2} \neq \mathrm{L}_{\alpha-1}, \mathrm{~L}_{\alpha}$. Suppose $M$ intersects $L_{\alpha-1}\left(L_{\alpha}\right)$ in a point $v$. If $v \neq x$ then apply (**) to $v$ and $x$. It follows that $M \in[v, x]$ intersects $K_{1}(y) \in[x, v]$ for $M$ misses $K_{2} \in[x, v]$. If $v x$ then $v=L_{\alpha-1} K_{i}\left(v=L_{\alpha} K_{i}\right)$ for some $i$. By Lemma 10 applied to $x$ and $L_{\alpha-1}\left(L_{\alpha}\right)$ it follows that $M$ intersects $K_{1}(y) \in\left[x, L_{\alpha-1}\right]\left(K_{1}(y) \in\left[x, L_{\alpha}\right]\right)$, for $M$ does not intersect $K_{2} \in\left[x, L_{\alpha-1}\right]\left(K_{2} \in\left[x, L_{\alpha}\right]\right)$.

COROLLARY. The line $\mathrm{K}_{1}(\mathrm{y})$ is the same for all $\mathrm{y} \in \mathrm{L}, \mathrm{y} \nmid \mathrm{x}$.
LEMMA 12. Let $\mathrm{x} \in \mathrm{X}, \mathrm{L} \in \mathrm{L}$ such that x is coltinear with a points $\mathrm{z}_{2}, \mathrm{z}_{3}, \ldots, \mathrm{z}_{\alpha+1}$ on L , Fut $\mathrm{K}_{\mathrm{i}}=\mathrm{xz}_{\mathrm{i}}, \mathrm{i}=2, \ldots, \alpha+1$ and let $\mathrm{K}_{1}$ be defined by $\left\{\mathrm{K}_{1}, \mathrm{~K}_{2}, \ldots, \mathrm{~K}_{\alpha+1}\right\}=[\mathrm{x}, \mathrm{y}]$ for any $\mathrm{y} \in \mathrm{L}, \mathrm{y} \nmid \mathrm{x}$. Then every line which intersects $\mathrm{K}_{1}$ and $a \mathrm{~K}_{\mathrm{i}}(\mathrm{i} \neq 1)$ not in x , intersects L and therefore exactly a innes of $\left\{\mathrm{K}_{1}, \ldots, \mathrm{~K}_{\alpha+1}\right\}$.

PROOF. Fix $i \in\{2, \ldots, \alpha+1\}$. The number of pairs ( $u, v$ ) such that
$u \in K_{1} \backslash\{x\}, v \in K_{i} \backslash\{x\}$, unv equals $s(\alpha-1)$. If $y \in L, y \notin x$ and $[y, x]=$ $\left\{L_{1}=\mathrm{L}, \mathrm{L}_{2}, \ldots, \mathrm{~L}_{\alpha+1}\right\}$, then each of the $\alpha-1$ 1ines $\mathrm{L}_{2}, \mathrm{~L}_{3}, \ldots, \mathrm{~L}_{\mathrm{i}-1}, \mathrm{~L}_{\mathrm{i}+1}, \ldots, \mathrm{~L}_{\alpha+1}$ gives rise to such a pair ( $u, v$ ). Each point $z_{j}, j=2,3, \ldots, i-1, i+1, \ldots, \alpha+1$ is on $\alpha-1$ lines $\neq K_{j}, \mathrm{~L}$ which intersect $\alpha$ lines of $\left\{\mathrm{K}_{1}, \ldots, \mathrm{~K}_{\alpha+1}\right\}$. They all intersect $K_{1}$ by Lemma 11 and no two miss the same $K_{k}$ since otherwise some $K_{\ell}$ would be hit $\alpha+1$ times. Thus each point $z_{j}, j=2,3, \ldots, i-1, i+1, \ldots, \alpha+1$ gives rise to ( $\alpha-2$ ) pairs ( $u, v$ ). Finally there are ( $\alpha-1$ ) pairs ( $u, v$ ) with $v=z_{i}$. In all, the lines intersecting $L$ contain $(s+1-\alpha)(\alpha-1)+$ $(\alpha-1)(\alpha-2)+(\alpha-1)=s(\alpha-1)$, i.e. a11, pairs (u,v). $\square$

If in Lemma 12 we replace $L=L_{1}$ by a 1 ine $L_{j}$ missing $K_{j}$, then it follows that every line intersecting two lines of $\left\{\mathrm{K}_{1}, \ldots, \mathrm{~K}_{\alpha+1}\right\}$ not in x , intersects exactly $\alpha$ lines of $\left\{\mathrm{K}_{1}, \ldots, \mathrm{~K}_{\alpha+1}\right\}$. Using this result and the foregoing lemmas we can now proceed as in the case $\mu=\alpha^{2}$. For any two intersecting lines $L_{1}, L_{2}$ we can define in an obvious way a partial geometry $\left\langle\mathrm{L}_{1}, \mathrm{~L}_{2}\right\rangle=\left(\mathrm{X}\left(\mathrm{L}_{1}, \mathrm{~L}_{2}\right), L\left(\mathrm{~L}_{1}, \mathrm{~L}_{2}\right)\right)$, now with parameters $\tilde{\mathrm{s}}=\mathrm{s}, \tilde{\mathrm{t}}=\alpha, \tilde{\alpha}=\alpha$ (so $\left\langle L_{1}, L_{2}>\right.$ is an $(\alpha+1)$-net of order $\left.s+1\right)$. Again $\tilde{\mu}=\tilde{\alpha}(\tilde{t}+1)=\alpha(\alpha+1)=\mu$, so with the same proof as the proof of fheorem 4 we have the following theorem.

THEOREM 8. Let ( $\mathrm{X}, \mathrm{L}$ ) be a s.p.g. with parameters $\mathrm{s}, \mathrm{t}, \alpha, \mu=\alpha(\alpha+1)$. If $\alpha \geq 4$, $s \geq f(\alpha)$ ( $f$ as in Lemma 9) antl $t \geq \alpha+1$ (i.e. if ( $\mathrm{X}, \mathrm{L}$ ) is not a p.g.), then (X,L) satisfies (D).

Fix a (d-2)-dimensional subspace $S$ of $P G(d, q), q$ a prime power, $d \in N$. Then with the lines of $P G(d, q)$ which have no point with $S$ in common as "points" and with the planes of PG(d,q) intersecting $S$ in exactly one point as "lines" and with the natural incidence relation, one obtains a s.p.g. with parameters $s=q^{2}-1, t=(q-1)^{-1}\left(q^{d-1}-1\right)-1, a=q, \mu=q(q+1)$.

DEBROEY [1] showed that a s.p.g. with parameters $s, t, \alpha \geq 2, \mu=$ $=\alpha(\alpha+1)$ and satisfying ( $D$ ) is of this type. Combining this result with Theorems 7 and 8 we arrive at the following theorem.

THEOREM 9. Let ( $\mathrm{X}, \mathrm{L}$ ) be a s.p.g. with parameters $\mathrm{s}, \mathrm{t}, \alpha, \mu=\alpha(\alpha+1)$ which is not a p.g.. If $\alpha=2$ and not $s=t=28$ or if $\alpha \geq 4$ and $s \geq f(\alpha)$, then $(\mathrm{X}, \mathrm{L})$ is isomorphic to a s.p.g. consisting of the lines in PG( $\mathrm{d}, \mathrm{q})$ missing a given (d-2)-dimensional subspace of $\mathrm{PG}(\mathrm{d}, \mathrm{q})$ and the planes inter-
secting this subspace in one point. In particular $s=q^{2}-1$, $\mathrm{t}=(\mathrm{q}-1)^{-1}\left(\mathrm{q}^{\mathrm{d}-1}-1\right)-1, \alpha=\mathrm{q}, \mathrm{p}=\mathrm{q}(\mathrm{q}+1)$ for some prime power q and $\mathrm{d} \in \mathbf{N}$ and any s.p.g. with these parameters with $q \neq 3$ and $\mathrm{d} \geq 4$ is of this type. REFERENCES
[1] DEBROEY, I., Semi partial geometries satisfying the diagonal axiom, J. Geometry, vol 13/2 (1979) 171-190.
[2] Debroey, I. \& J.A. THAS, on semi partial geometries, J. Comb. Th.
(A) 25 (1978) 242-250.

## SAMENVATTING

Dit proefschrift bestaat uit vijf artikelen en een inleidend hoofdstuk. In elk van de vijf artikelen wordt een karakterisering gegeven van een object uit de eindige meetkunde. Het inleidende hoofdstuk bestaat uit een overzicht van vergelijkbare resultaten uit de literatuur (de stelling van Veblen \& Young over projectieve ruimtes, de stelling van ostrom \& Wagner over projectieve vlakken met een 2-transitieve automorfismengroep, de stelling van Buekenhout \& Shult over polaixe ruimtes, etc.), en inleidingen in elk van de vijf artikelen.

Het eerste artikel gaat over bijna affiene vlakken. Evenals bij gewone affiene vlakken is het ook hier mogelijk het begrip translatie te definiexen. Aangetoond wordt dat het bestaan van translaties equivalent is met de geldigheid van een "Stelling van Desargues", en dat bijna affiene vlakken met een transitieve groep van translaties op een bepaalde algebraische manier kunnen worden beschreven.

In het tweede artikel wordt aangetoond dat er een verband bestaat tussen bijna affiene vlakken en Minkowski vlakken. Dit gegeven wordt gebruikt om een meetkundige karakterisering te geven van alle, tot nu toe bekende, Minkowski vlakken. In essentie komt deze karakterisering neer op de eis dat alle bijna affiene vlakken die met een Minkowski vlak zijn geassocieerd, moeten voldoen aan de stelling van Desargues.

In het derde artikel wordt een tweede karakterisering gegeven van de op dit moment bekende Minkowski vlakken. Het blijken precies die Minkowski vlakken te zijn warvan de automorfismengroep transitief is op paren nietcollineaire punten.

Het vierde artikel qeeft een meetkundige karakterisering van de klassieke unital (dit is het $2-\left(q^{3}+1, q+1,1\right)$ design van de absolute punten en niet absolute lijnen van een unitaire polariteit van $\left.P G\left(2, q^{2}\right)\right)$. De gekozen meetkundige condities zijn zodanig dat een op de punten 2 -transitieve groep van automorfismen geconstrueerd kan worden, die vervolgens gefdentificeerd wordt als $\operatorname{PSU}\left(3, q^{2}\right)$.

Het vijfde en laatste artikel geeft een karakterisering van twee klassen van semi-partiele meetkundes die geconstrueerd kunnen worden uit projectieve ruimtes. Bij deze kaxakterisering wordt alleen uitgegaan van de speciale vorm van de parameters. Het doel wordt hier bereikt door aan te tonen dat in deze semi-partiele meetkundes de duale versie van het axioma van Pasch geldt.

## STELLINGEN

behorende bij het proefschrift
van
H.A. WILBRINK

1. De bewijstechniek van [1] is mogelijkerwijs ook te gebruiken om het niet bestaan van een sterk reguliere graaf op 99 punten van graad 14 aan te tonen.
[1] H.A. Wilbrink \& A.E. Brouwer, A (57,14,1) strongly regular graph does not exist. proc. KNAW A 86 (1) , 1983.
2. Ex bestaat tenminste een (symmetrisch) $2-(49,16,5)$ design.
A.E. Brouwer \& H.A. Wilbrink, A symmetric design with parameters $2-(49,16,5)$, to appear.
3. De punten en lijnen die geheel buiten een niet ontaarde hyperkwadriek in $P G(2 n-1,2)$ liggen, vormen een semi-partiele meetkunde.
4. Veronderstel dat een rang 3 Zara graaf de volgende eigenschap heeft. Voor ieder tweetal disjuncte vlakken bestaan er partities in lijnen van die vlakken zo dat iedere lijn van elke partitie in een vlak is met een lijn uit de andere partitie. Dan vormen de vlakken en lijnen met de natuurlijke incidentie de punten en lijnen van een bijna-zeshoek.
A. Blokhuis, Few-distance sets, Proefschrift T.H.E., 1983.
E. Shult \& A. Yanushka, Near n-gons \& line systems, Geom. Dedicata 9 (1980), 1-72.
5. Vermoedelijk geldt de volgende stelling. Als $n$ de orde is van een projectief vlak met een reguliere abelse automorfismengroep en $p$ is een priemdelex van $n$, dan is $n=p$ of $p^{2}$ deelt $n$. Voor $p=2$ en $p=3$ is dit bewezen.
H.A. Wilbrink, A note on planar difference sets, to appear.
6. De grafen op de inwendige en uitwendige punten van een niet ontaarde hyperkwadriek in PG(2n,5), met als kanten de paren onderling loodrechte punten, zijn sterk regulier.
7. Veel bewijzen in de combinatoxiek kunnen met $50 \%$ worden ingekort door gebruik te maken van matrices. Vergelijk: E. Artin, Geometrio Algebra, Interscience, New York, 1957, p 14.
8. Het adagium "een plaatje is geen bewijs" dient zeker niet te worden geinterpreteerd als een aanbeveling tot het niet gebruiken van plaatjes in de wiskunde.
9. Het gebruik van computers binnen de wiskunde kan een remmende invloed hebben op de ontwikkeling van de wiskunde, en dient daarom met de nodige terughoudendheid te gebeuren.

[^0]:    1. Benz, W.: Vorlesungen über Geometrie der Algebren, Springer-Verlag, Berlin, New York, 1973.
    2. Bruen, A. and Levinger, B.: 'A Theorem on Permutations of a Finite Field’, Can. J. Math. 25 (1973), 1060-1065.
    3. Dembowski, P.: Finife Geometries, Springer-Verlag, Berlin, Heidelberg, New York, 1968.
    4. Feilt, W.: 'On a Class of Doubly Transitive Permutation Groups, III. J. Math. 4 (1960), 170-186.
    5. Heise, W. and Karzel, H.: 'Symmetrische Minkowski-Ebenen', J. Geometry 3 (1973), 5-20.
    6. Huppert, B. : 'Scharf Dreifach Transitive Permutationsgruppen', Arch. Math. 13 (1962), 61-72.
