# Normal forms for real surfaces in $\mathbf{C}^{2}$ near complex tangents and hyperbolic surface transformations 

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## 0. Introduction

It is well known that the complex analytical properties of a real submanifold $M$ in the complex space $\mathbf{C}^{n}$ are most accessible through consideration of the complex tangents to $M$. The properties we have in mind are related to the behavior of holomorphic functions on or near $M$ and to the behavior of $M$ under biholomorphic transformation. The case in which $M$ is a real hypersurface is most familiar, while much less is known for higher codimension. In this paper we consider the critical case of a real $n$ dimensional manifold $M$ in $C^{n}$, which we also assume to be real analytic. At a generic point $M$ is locally equivalent to the standard $\mathbf{R}^{n}$ in $\mathbf{C}^{n}$. However, near a complex tangent $M$ may aquire a non-trivial local hull of holomorphy and other biholomorphic invariants.

We begin with the simplest non-trivial case, which is a surface $M^{2} \subset C^{2}$ with an isolated, suitably non-degenerate complex tangent. Here one already encounters a rich structure and non-trivial problems. In coordinates $z_{j}=x_{j}+i y_{j}, j=1,2, M$ may be written locally as

$$
\begin{aligned}
& R(z, \bar{z})=-z_{2}+q\left(z_{1}, \bar{z}_{1}\right)+\ldots=0 \\
& q=\gamma z_{1}^{2}+z_{1} \bar{z}_{1}+\gamma \bar{z}_{1}^{2}, \quad 0 \leqslant \gamma<\infty
\end{aligned}
$$

The $z_{1}$-axis is tangent to $M$ at the origin. $M$, or more precisely, this complex tangent is said to be elliptic if $0 \leqslant \gamma<1 / 2$, hyperbolic if $1 / 2<\gamma$, or parabolic if $\gamma=1 / 2$. We shall prove the following theorem.

[^0]THEOREM 1. Let $M$ be a real analytic surface in $\mathbf{C}^{2}$ with an elliptic complex tangent at a point $p$ with $0<\gamma<1 / 2$. Then there exists a holomorphic coordinate system $\left(z_{1}, z_{2}\right)$ in which $p=0$, and $M$ has locally the form

$$
x_{2}=z_{1} \bar{z}_{1}+\Gamma\left(x_{2}\right)\left(z_{1}^{2}+\bar{z}_{1}^{2}\right), \quad y_{2}=0,
$$

where $\Gamma=\gamma+\delta x_{2}^{s}, \delta= \pm 1, s \in \mathbf{Z}^{+}$, or $\Gamma=\gamma(s=\infty)$. The quantities $\gamma, \delta$, sform a complete system of biholomorphic invariants for $M$ near $p$.

A consequence of Theorem 1 is that the local hull of holomorphy of $M$ near $p$ is precisely the real analytic 3 -manifold-with-boundary $\tilde{M}: x_{2} \geqslant z_{1} \bar{z}_{1}+\Gamma\left(x_{2}\right)\left(z_{1}^{2}+\bar{z}_{1}^{2}\right), y_{2}=0$. $\tilde{M}$ is the union of a one-parameter family of ellipses, the boundaries of which are the curves on $M$ gotten by setting $x_{2}=c>0$. Another consequence is that such an $M$ always admits the biholomorphic involution corresponding to $\left(z_{1}, z_{2}\right) \mapsto\left(-z_{1}, z_{2}\right)$. It is interesting to note that $M$ is locally equivalent to an algebraic surface.

We also have the analogue of Theorem 1 in the $n$-dimensional case. This theorem will be reduced to a seemingly unrelated problem, namely that of a normal form for a pair of involutions $\tau_{1}, \tau_{2}$ which are holomorphic mappings in a neighborhood of a common fixed point $p$ in $\mathbf{C}^{2}$. They are subjected to biholomorphic mappings $\psi$ keeping $p$ fixed, by replacing $\tau_{j}$ by $\psi^{-1} \tau_{j} \psi$. We ask for a classification of the pairs of $\left(\tau_{1}, \tau_{2}\right)$, and more generally of the group generated by the $\tau_{j}$, under the pseudo group of biholomorphic mappings near $p$.

Taking $\xi, \eta$ as coordinates in $\mathbf{C}^{2}, p$ as the origin and the linearized maps $\left.d \tau_{j}\right|_{0}$ as

$$
d \tau_{j}:(\xi, \eta) \mapsto\left(\lambda_{j} \eta, \lambda_{j}^{-1} \xi\right) \quad \text { with } \quad \lambda_{1}=\lambda_{2}^{-1}=\lambda \neq 0
$$

we can state our result as follows:
THEOREM 2. If $|\lambda| \neq 1$ then there exists a biholomorphic mapping $\psi$ near the origin with $\psi(0)=0$, taking the two given holomorphic involutions $\tau_{j}$ into

$$
\psi^{-1} \tau_{j} \psi:(\xi, \eta) \mapsto\left(\Lambda_{j}(\xi \eta) \eta, \Lambda_{j}(\xi \eta)^{-1} \xi\right)
$$

where

$$
\Lambda_{1}=\Lambda_{2}^{-1}=\lambda+\delta(\xi \eta)^{s}, \quad \delta=1,0, s \geqslant 1
$$

For our application we will have to consider these holomorphic involutions $\tau_{j}$ in conjunction with an antiholomorphic involution $\varrho$ describing the reality condition and satisfying

$$
\tau_{1} \varrho=\varrho \tau_{2} .
$$

This leads to a finer classification of ( $\tau_{1}, \tau_{2}$ ) under the group of biholomorphic mappings $\psi$ which commute with $\varrho$. In particular, we will have $\lambda$ real or $|\lambda|=1$ in correspondence to the elliptic or hyperbolic quadric.

In Theorem 2 we had to rule out that $\lambda$ lies on the unit circle. Actually, if $|\lambda|=1, \lambda$ not a root of unity, one can still find a formal power series expansion for $\psi$, but in general one has to expect divergence of these series. This is a "small divisor problem" as one encounters it in Celestial Mechanics. The product $\varphi=\tau_{1} \tau_{2}$ is the crucial map which has to be normalized. Its linearization at the origin $\left.d \varphi\right|_{0}$ has the eigenvalues $\lambda^{2}, \lambda^{-2}$.

However, in celestial mechanics one restricts oneself to area-preserving mappings, and the analogous equivalence problem was studied by G. D. Birkhoff [3]. In our case the area-preserving property is replaced by the condition

$$
\varphi^{-1}=\tau_{2} \tau_{1}=\tau_{1}^{-1} \varphi \tau_{1}
$$

which corresponds to "reversible" systems of differential equations. Mappings of this nature which can be represented as a product of involutions actually also played a role in Birkhoff's study of the restricted three body problems [2].

In case $\lambda$ is not on the unit circle one has no difficulty of small divisors but the corresponding convergence proofs are not straightforward. We apply here a refinement of the majorant method as it was developed for area-preserving mappings in [12] and [11].

One may be led to the involutions $\tau_{1}, \tau_{2}, \varrho$ of Theorem 2 by the problem of characterizing intrinsically the trace $f$ on $M$ of a function $g$ holomorphic in a neighborhood of $M$. The key to this is to complexify $M$. Replacing $\bar{z}$ by independent complex variables $w$ in the equation for $M$ gives the complex analytic surface

$$
\mathfrak{M}=\left\{(z, w) \in \mathbf{C}^{4}: R(z, w)=0, \bar{R}(w, z)=0\right\}
$$

If the natural projections $\pi_{1}(z, w)=z, \pi_{2}(z, w)=w$ are restricted to $\mathfrak{M}$, then $f$ and $g$ are related by $f=g \circ \pi_{1}$. For $\gamma \neq 0, \pi_{1}$ and $\pi_{2}$ are two-fold branched coverings. The covering transformations $\tau_{2}, \tau_{1}$ are holomorphic involutions on $\mathfrak{M}$ fixing the origin. The condition $f \circ \tau_{2}=f$ is an intrinsic characterization of the restriction of a function holomorphic in $z$. It is a discrete analogue of the local characterization [9] by H . Lewy of the restriction of a holomorphic function to a strongly pseudo-convex real hypersurface. $\tau_{2}$ corresponds to the tangential Cauchy-Riemann operator, and the mapping $\varphi$ is a discrete version of the Levi-form. In the elliptic case $\varphi$ can be embedded in a flow $\varphi^{t}, t$
complex. The orbits of this flow intersect $M$ precisely in the curves bounding the above mentioned analytic discs.

If the surface $M$ is elliptic, then $\lambda$ is real, $\lambda \neq \pm 1$, and the origin is a hyperbolic fixed point of $\varphi$ in the sense of mappings. In this case we have a satisfactory theory. If $M$ is hyperbolic, $\varphi$ is elliptic. The subtleties of the theory of elliptic mappings, e.g. small divisors, make the theory of hyperbolic surfaces much more difficult.

Previously it was known to Bishop [4] that $\gamma$ is a biholomorphic invariant. He also proved in the elliptic case the existence of a one-parameter family of analytic discs with boundaries on $M^{2} \subset \mathbf{C}^{2}$ near the complex tangent $p$. Hunt [7] investigated further the regularity of $\tilde{M}$, the union of these discs. In [8] it was shown that $\bar{M}$ is a $C^{\infty}$ manifold-with-boundary for $0 \leqslant \gamma<1 / 2$, and that the discs are unique. In [1] Bedford and Gaveau consider hulls of holomorphy from a global viewpoint.

In section one of this paper we discuss the connection between surfaces and involutions. In fact, we show the equivalence of certain complex surfaces $\mathfrak{M}$ with suitable pairs of involutions $\tau_{1}, \tau_{2}$. We dicuss thoroughly in section 2 the quadric surfaces, which correspond to pairs of linear involutions. Here the basic phenomena are clearly revealed. Section 3 deals with pairs of non-linear involutions on a formal level, and section 4 contains the convergence proof for Theorem 2 . In section 5 these results are applied to derive the normal form for the manifold $M$.

In section 6 we discuss hyperbolic surfaces. In particular, we show divergence of the transformation into normal form for an example, using ideas previously developed for area-preserving Cremona transformations [10].

## 1. Surfaces and involutions

Let $M$ be a smooth real analytic surface in $\mathbf{C}^{2}$. It may be described locally by two independent real equations or by one complex equation,

$$
\begin{equation*}
M: R(z, \bar{z})=0, \quad \dot{R}(\bar{z}, z)=0, \quad d R \wedge d \bar{R} \neq 0, \tag{1.1}
\end{equation*}
$$

where $R$ is a power series in $z=\left(z_{1}, z_{2}\right)$ and $\bar{z}$. We wish to investigate the local properties of $M$ under the pseudo-group of local biholomorphic transformations

$$
z^{\prime}=f(z), \quad z^{\prime}=f(\bar{z}) .
$$

We assume that the point $z=0$ lies on $M$. By interchanging $R$ and $\dot{R}$ we may assume that the holomorphic linear term in $R$ is non-zero. After introducing this linear function as a
new $z_{2}$ variable, we may assume that $M$ has the form $z_{2}=p \bar{z}_{1}+q \bar{z}_{2}+\ldots$ and after a further transformation we achieve $q=0$, so that

$$
\begin{equation*}
z_{2}=p \bar{z}_{1}+F\left(z_{1} \bar{z}_{1}\right), \quad F=O\left(\left|z_{1}\right|^{2}\right) \tag{1.2}
\end{equation*}
$$

If $p \neq 0$ we may introduce new coordinates $z^{\prime}$ by $z_{2}=p z_{2}^{\prime}+F\left(z_{1}^{\prime}, z_{2}^{\prime}\right), z_{1}=z_{1}^{\prime} . M$ then goes over into the totally real plane $z_{2}^{\prime}=\bar{z}_{1}^{\prime}$. Hence, $M$ has no invariants near such a point. We henceforth assume that $p=0$, so that the $z_{1}$-axis is a complex tangent to $M$ at the origin. $M$ is then given by $z_{2}=a z_{1}^{2}+b z_{1} \bar{z}_{1}+c \bar{z}_{1}^{2}+\ldots$. We make the non-degeneracy assumption that $b \neq 0$. Then by a quadratic change of $z_{2}$ we may achieve $b=1, a=\bar{c}=\gamma$. A rotation of $z_{1}$ makes $0 \leqslant \gamma$. The surface $M$ is now assumed to have the form

$$
M: \begin{align*}
& z_{2}=F\left(z_{1}, \bar{z}_{1}\right)=q\left(z_{1}, \bar{z}_{1}\right)+H\left(z_{1}, \bar{z}_{1}\right), \\
& q=\gamma z_{1}^{2}+z_{1} \bar{z}_{1}+\gamma \bar{z}_{1}^{2}, \quad 0 \leqslant \gamma<\infty, H=O\left(\left|z_{1}\right|^{3}\right) . \tag{1.3}
\end{align*}
$$

The non-negative number $\gamma$ is a biholomorphic invariant of $M$ first considered by Bishop [4]. The complex tangent is elliptic if $0 \leqslant \gamma<1 / 2$, parabolic if $\gamma=1 / 2$, or hyperbolic if $1 / 2<\gamma<\infty$.

For our investigation it will be necessary to characterize those real analytic functions on the surface $M$ which are the restrictions of functions holomorphic in some neighborhood of $M$. This is facilitated by complexifying $M$. We replace $\bar{z}$ by independent variables $w=\left(w_{1}, w_{2}\right)$ in (1.1) and define a smooth complex analytic surface $\mathfrak{M}$ in $C^{4}$ by

$$
\mathfrak{M}: R(z, w)=0, \bar{R}(w, z)=0
$$

Complex conjugation $(z, \bar{z}) \mapsto(\bar{z}, z)$ goes over into the anti-holomorphic involution

$$
\varrho(z, w)=(\bar{w}, \bar{z})
$$

More generally we consider a complex surface

$$
\mathfrak{M}: R(z, w)=0, \quad S(z, w)=0, \quad d R \wedge d S \neq 0
$$

passing through the origin of $\mathbf{C}^{4}$ under the wider group of transformations

$$
\begin{equation*}
z^{\prime}=f(z), \quad w^{\prime}=g(w) \tag{1.4}
\end{equation*}
$$

Such an $\mathfrak{M}$ comes from a real surface $M \subset C^{2}$ if and only if $\varrho \mathbb{M}=\mathbb{M}$, and such a transformation is induced by a holomorphic mapping of $\mathbf{C}^{2}$ if and only if $f(z)=\bar{g}(z)$. (The bar indicates complex conjugation of the coefficients only in the series $g$.)

There are two invariant projections on $\mathbf{C}^{4}, \pi_{1}(z, w)=z, \pi_{2}(z, w)=w$. They are related by $\pi_{2}=c \pi_{1} \varrho$, where $c$ denotes complex conjugation. We denote the restrictions of $\pi_{1}$ and $\pi_{2}$ to $\mathfrak{M}$ by the same symbols. The real and imaginary parts of $z_{1}$ may be taken as coordinates on $M$, and $\left(z_{1}, w_{1}\right)$ as coordinates on $\mathfrak{M}$. A real analytic function $f=f\left(z_{1}, z_{1}\right)$ on $M$ may be continued locally to a function $f=f\left(z_{1}, w_{1}\right)$ holomorphic on $\mathfrak{M}$. The original function is the restriction of a holomorphic function if and only if the extended function $f$ satisfies $f=g \circ \pi_{1}$ for some function $g=g(z)$ holomorphic in $z$. Similarly $f$ is the restriction of an anti-holomorphic function if and only if the extended $f$ satisfies $f=g \circ \pi_{2}, g=g(w) . f$ is real if and only if $f \circ \varrho=c f$.

The possible linear structure of $\mathfrak{M}$ is more varied. To describe it let $P_{z}=\{w=0\}$ and $P_{w}=\{z=0\}$ denote the $z$ and $w$ coordinate planes, and $P$ denote the tangent complex two-plane to $\mathfrak{M}$ at the origin. There are four possibilities;
(1) $P$ is totally real: $\operatorname{dim} P \cap P_{z}=\operatorname{dim} P \cap P_{w}=0$,
(2) $P$ is partially holomorphic: $\operatorname{dim} P \cap P_{z} \geqslant 1$,
(3) $P$ is partially anti-holomorphic: $\operatorname{dim} P \cap P_{w} \geqslant 1$,
(4) $P$ is complex: $\operatorname{dim} P \cap P_{z}=\operatorname{dim} P \cap P_{w}=1$.

We shall study $\mathfrak{M}$ only in a neighborhood of a point at which its tangent plane $P$ is complex (type (4)). Generically through such a point there exist a curve $C_{1}$ of points at which the tangent plane is of type (2) and a curve $C_{2}$ of points at which it is of type (3). Locally $\mathfrak{M}$ is given by

$$
\begin{align*}
& z_{2}=F\left(z_{1}, w_{1}\right)=(q+H)\left(z_{1}, w_{1}\right)  \tag{1.5}\\
& w_{2}=G\left(z_{1}, w_{1}\right)=(p+K)\left(z_{1}, w_{1}\right)
\end{align*}
$$

The quadratic terms $q$ and $p$ are both assumed to have a non-zero $z_{1} w_{1}$-term, and so may be put into the form

$$
q=a z_{1}^{2}+z_{1} w_{1}+a w_{1}^{2}, \quad p=b z_{1}^{2}+z_{1} w_{1}+b w_{1}^{2}
$$

via a transformation (1.4). The product $a b$ is invariant under (1.4). If $a b \neq 0$ then by a substitution $\left(z_{1}, w_{1}\right) \mapsto\left(\alpha z_{1}, \alpha^{-1} w_{1}\right)$ we may achieve $a=b=\gamma \in \mathbf{C} . \gamma$ is then invariant up to sign.

We now assume that

$$
\begin{equation*}
p=q=\gamma z_{1}^{2}+z_{1} w_{1}+\gamma w_{1}^{2}, \quad \gamma \neq 0 . \tag{1.6}
\end{equation*}
$$

In this case, when restricted to $\mathfrak{M}$, the projections

$$
\pi_{1}\left(z_{1}, w_{1}\right)=\left(z_{1}, F\left(z_{1}, w_{1}\right)\right), \quad \pi_{2}\left(z_{1}, w_{1}\right)=\left(w_{1}, G\left(z_{1}, w_{1}\right)\right)
$$

are locally two-fold branched coverings. The branch locus of $\pi_{1}$ is given by

$$
z_{2}=F\left(z_{1}, w_{1}\right), \quad F_{w_{1}}\left(z_{1}, w_{1}\right)=0
$$

which is a smooth curve in the $\left(z_{1}, z_{2}\right)$-plane since $F_{w_{1} w_{1}}(0, \overline{0})=2 \gamma \neq 0$. Likewise, the branch locus of $\pi_{2}$ is a non-singular curve in the $w$-plane. The equation $w=w^{\prime}$ or

$$
q\left(z_{1}^{\prime}, w_{1}\right)-q\left(z_{1}, w_{1}\right)=K\left(z_{1}, w_{1}\right)-K\left(z_{1}^{\prime}, w_{1}\right)
$$

together with (1.5) generally have a unique solution $z^{\prime} \neq z$. By the implicit function theorem they define a local self-mapping of $\mathfrak{M}$

$$
\begin{align*}
& \tau_{1}: \begin{array}{l}
z_{1}^{\prime}=-z_{1}-\frac{1}{\gamma} w_{1}+h_{1}\left(z_{1}, w_{1}\right) \\
w_{1}^{\prime}=w_{1}
\end{array}, \tag{1.7}
\end{align*}
$$

which is an involution, $\tau_{1}^{2}=$ id. Similarly $z^{\prime}=z$ or $F\left(z_{1}, w_{1}^{\prime}\right)=F\left(z_{1}, w_{1}\right)$ and (1.5) determines a second involution

$$
\begin{align*}
& z_{1}^{\prime}=z_{1} \\
& \tau_{2}:  \tag{1.8}\\
& w_{1}^{\prime}=-\frac{1}{\gamma} z_{1}-w_{1}+h_{2}\left(z_{1}, w_{1}\right)
\end{align*}
$$

$\tau_{1}$ and $\tau_{2}$ are the covering transformations for $\pi_{2}$ and $\pi_{1}, \pi_{1} \tau_{2}=\pi_{1}, \pi_{2} \tau_{1}=\pi_{2}$. The fixed point sets of $\tau_{1}$ and $\tau_{2}$ are the curves $C_{1}$ and $C_{2}$ mentioned above. If $\mathfrak{M}$ satisfies the reality condition, then $\tau_{1} \varrho=\varrho \tau_{2}$.

We may now characterize the trace $f$ on $\mathfrak{M}$ of a function $g(z)$ holomorphic in $z$. Since $f=g \circ \pi_{1}$, we necessarily have $f \circ \tau_{2}=f$. Conversely, suppose $f=f\left(z_{1}, w_{1}\right)$ is analytic in $\left(z_{1}, w_{1}\right)$ and invariant under $\tau_{2}$. There then exists a single-valued function $g=g\left(z_{1}, z_{2}\right)$, defined and holomorphic on the base of the branched covering $\pi_{1}$ away from its branch locus, satisfying $f=g \circ \pi_{1}$, or $f\left(z_{1}, w_{1}\right)=g\left(z_{1}, F\left(z_{1}, w_{1}\right)\right)$. Since $g$ is bounded it extends to be holomorphic in a neighborhood of $z=0$ by the Riemann extension theorem. We may also say that the functions $z_{1}$ and $F\left(z_{1}, w_{1}\right)$ generate the algebra of $\tau_{2}$-invariants. The condition $f \circ \tau_{1}=f$ characterizes the trace on $\mathfrak{M}$ of a function holomorphic in $w$ (anti-holomorphic in a neighborhood of $M$ ).

The mapping $\varphi=\tau_{1} \tau_{2}$,

$$
\varphi: \begin{align*}
& z_{1}^{\prime}=-\left(1-\gamma^{-2}\right) z_{1}+\gamma^{-1} w_{1}+\ldots  \tag{1.9}\\
& w_{1}^{\prime}=-\gamma^{-1} z_{1}-w_{1}+\ldots
\end{align*}
$$

is also very important for the study of the surface $M$. It has the origin as an isolated fixed point. We shall see in section 2 that this fixed point is hyperbolic if $M$ is elliptic and elliptic if $M$ is hyperbolic.

The condition $f \circ \tau_{2}=f$ is an analogue of the tangential Cauchy-Riemann equations on a real hypersurface in $\mathbf{C}^{n}$. In fact, such an analytic $M^{2 n-1} \subset \mathbf{C}^{n}$ yields, upon complexification, $\mathfrak{M}^{2 n-1} \subset \mathbf{C}^{2 n}$. The two projections $\pi_{1}, \pi_{2}: \mathfrak{M}^{2 n-1} \rightarrow \mathbf{C}^{n}$ each have rank $n$ and $(n-1)$-dimensional fibers. The tangents to the fibers of $\pi_{1}$ are spanned by the $n-1$ independent complexified tangential Cauchy-Riemann operators. A function on $\mathfrak{M}^{2 n-1}$ annihilated by these operators is constant on the fibers of $\pi_{1}$ and so comes from a function holomorphic in $z$ alone. In the case $n=2$, there is only one independent tangential vector field $P$ of type ( 1,0 ). By complexification, $(P, \bar{P}$ ) goes over to ( $P, Q$ ) on $\mathfrak{M}^{3}$. We consider the flows

$$
\psi_{1}^{t}=\exp (t P), \quad \psi_{2}^{t}=\exp (t \mathrm{Q})
$$

with $t$ complex. They commute when $[P, Q]=0$, which implies that the Levi-form of $M$ vanishes. Thus the commutator $\varphi^{2}=\tau_{1} \tau_{2} \tau_{1}^{-1} \tau_{2}^{-1}$ may be thought of as a discrete analogue of the Levi-form. Under the assumption that the linear part of $\varphi$ is not nilpotent, we shall derive a normal form for $M$ in section 3 . This may be compared to the normal form in [5] for a non-degenerate real hypersurface.

To further emphasize the importance of $\tau_{1}$ and $\tau_{2}$, we next wish to show that two suitable such involutions defined and holomorphic in a neighborhood of and fixing the origin of $\mathbf{C}^{2}$ give rise to a surface $\mathfrak{M}$ in $\mathbf{C}^{4}$. Let the coordinates of $\mathbf{C}^{2}$ be denoted by $X=(x, y)^{t}$, and suppose

$$
\begin{align*}
& \tau_{j}: X^{\prime}=T_{j} X+h_{j}(X), \quad h_{j}=O\left(|X|^{2}\right),  \tag{1.10}\\
& T_{j}^{2}=1, \quad h_{j} \circ \tau_{j}=-T_{j} h_{j}, \quad j=1,2
\end{align*}
$$

For each $j$ we assume that the 2 by 2 matrix $T_{j}$ has a ( -1 )-eigenspace of dimension one, and consequently a one-dimensional ( +1 )-eigenspace. Let the eigenvectors be denoted by $v_{j}^{+}, v_{j}^{-}$. We further assume that each of the pairs of vectors $\left(v_{1}^{-}, v_{2}^{-}\right),\left(v_{1}^{-}, v_{2}^{+}\right),\left(v_{2}^{-}, v_{1}^{+}\right)$ is linearly independent. After a linear change of coordinates $(x, y)$, we may assume $v_{1}^{-}=(1,0)^{t}$ and $v_{2}^{-}=(0,1)^{t}$, and so

$$
T_{1}=\left(\begin{array}{rr}
-1 & c_{1}  \tag{1.11}\\
0 & 1
\end{array}\right), \quad T_{2}=\left(\begin{array}{rr}
1 & 0 \\
c_{2} & -1
\end{array}\right)
$$

The other two independence conditions are equivalent to $c_{1} \neq 0, c_{2} \neq 0$. More explicitly we now have

$$
\begin{align*}
& \tau_{1}: \begin{array}{l}
x^{\prime} \\
y^{\prime}
\end{array}=-x+c_{1}(x, y) \\
& x_{1}(x, y) \\
& x_{2}=x+g_{2}(x, y)  \tag{1.12}\\
& y^{\prime}=c_{2} x-y+f_{1}(x, y)
\end{align*}
$$

where the $f_{j}, g_{j}$ are of second order. We define four new functions by

$$
\begin{align*}
& z_{1}=y+y \circ \tau_{2}=c_{2} x+f_{2}(x, y) \\
& z_{2}=y \cdot y \circ \tau_{2}=c_{2} x y-y^{2}+y f_{2}(x, y) \\
& w_{1}=x+x \circ \tau_{1}=c_{1} y+f_{1}(x, y)  \tag{1.13}\\
& w_{2}=x \cdot x \circ \tau_{1}=c_{1} x y-x^{2}+x f_{1}(x, y)
\end{align*}
$$

Clearly, $(x, y) \mapsto(z, w)$ defines a map of rank two. If we use the first and third equation to eliminate $x$ and $y$, then the image is seen to lie on a surface

$$
\begin{aligned}
& z_{2}=c_{1}^{-1} z_{1} w_{1}-c_{1}^{-2} w_{1}^{2}+\ldots \\
& w_{2}=c_{2}^{-1} z_{1} w_{1}-c_{2}^{-2} z_{1}^{2}+\ldots
\end{aligned}
$$

This can be put into the form (1.5), (1.6) by a transformation (1.4). It is clear that $\tau_{1}$ and $\tau_{2}$, by fixing the functions $w$ and $z$, respectively, are the involutions induced by this embedding. Since ( $w_{1}, w_{2}$ ) and ( $z_{1}, z_{2}$ ) generate the functions invariant under $\tau_{1}$ and $\tau_{2}$, respectively, it follows that any other such regularly embedded surface realizing the $\tau_{j}$ is equivalent to this one via a transformation (1.4). We have proved the following.

PROPOSITION 1.1. Every analytic surface $(1.5,1.6)$ gives rise to an intrinsic pair of involutions (1.12). Conversely, every such pair (1.12) are the intrinsic involutions of some surface $(1.5,1.6)$ in $\mathbf{C}^{4}$.

Note that the two $(+1)$-eigenvectors $\left(v_{1}^{+}, v_{2}^{+}\right)$are dependent precisely when $c_{1} c_{2}=4$. This is the parabolic case, corresponding to $\gamma= \pm 1 / 2$ in (1.6).

An anti-holomorphic involution $\varrho$ fixing the origin in $\mathbf{C}^{\mathbf{2}}$ has the form

$$
\begin{gather*}
\varrho: X^{\prime}=P \bar{X}+k(\bar{X}), \quad k=O\left(|X|^{2}\right) \\
P \bar{P}=I, \quad k(\overline{\varrho(X)})=-P \bar{k}(X), \tag{1.14}
\end{gather*}
$$

the last two equations being equivalent to $\varrho^{2}=\mathrm{id}$. The fixed points of $\varrho$ are the solutions to the system $X=\varrho(X)$, which in view of the conditions (1.14), is equivalent to a single equation (1.2) with $p \neq 0$. If this surface is transformed into $\mathbf{R}^{2} \subset \mathbf{C}^{2}$, then $\varrho$ may be put into the form $\varrho(x, y)=(\bar{x}, \bar{y})$. The change $(x, y) \mapsto(x+i y, x-i y)$ gives the form $\varrho(x, y)=(\bar{y}, \bar{x})$. If the involutions (1.10) also satisfy $\tau_{1} \varrho=\varrho \tau_{2}$ for $\varrho$ in this form, then $z \circ \varrho=\bar{w}$, and (1.12) gives rise to a real surface (1.3) in $\mathbf{C}^{2}$.

Next, we consider a real analytic $n$-dimensional manifold in $\mathbf{C}^{n}$. Generically, such a manifold $M$ is totally real and so locally equivalent to the standard $\mathbf{R}^{n}$ in $\mathbf{C}^{n}$. We shall, however, study $M$ near a point $p$ at which it has a complex one-dimensional holomorphic tangent space. We use the following coordinate notation

$$
\begin{align*}
& z=\left(z_{1}, z_{\alpha}, z_{n}\right), z_{\alpha}=x_{\alpha}+i y_{\alpha}, \quad 2 \leqslant \alpha \leqslant n-1 \\
& x=\left(x_{2}, \ldots, x_{n-1}\right) . \tag{1.15}
\end{align*}
$$

The Greek indices $\alpha, \beta, \sigma$ will generally have the range from 2 to $n-1$ throughout this paper. These coordinates are initially chosen so that $p$ is the origin, the $z_{1}$-axis is the holomorphic tangent space to $M$ at $p$, the ( $z_{1}, x$ )-space is the real tangent space $T_{p}$ to $M$ at 0 , and $z_{n}=0$ is the complex envelope, $T_{p}+i T_{p}$, of this real tangent space. We may then express $M$ locally as a graph

$$
\begin{align*}
& z_{n}=F\left(z_{1}, \bar{z}_{1}, x\right)  \tag{1.16}\\
& y_{a}=f_{a}\left(z_{1}, \bar{z}_{1}, x\right)=f_{a}\left(\bar{z}_{1}, z_{1}, x\right)
\end{align*}
$$

where $F, f_{\alpha}$ begin with quadratic terms. Those in $F$ have the form $q+q_{1}+q_{2}$,

$$
\begin{aligned}
q & =a z_{1}^{2}+b z_{1} \bar{z}_{1}+c \bar{z}_{1}^{2} \\
q_{1} & =\sum a_{a} x_{\alpha} z_{1}+b_{\alpha} x_{a} \bar{z}_{1} \\
q_{2} & =\sum c_{a \beta} x_{\alpha} x_{\beta}
\end{aligned}
$$

As in the two dimensional case we make the non-degeneracy assumption that $b \neq 0$, and even that $b=1, a=c=\gamma, 0 \leqslant \gamma<\infty$.

Further simplification of $F, f_{\alpha}$ is made as follows. To eliminate the second term in $q_{1}$ we make the change $z_{1} \mapsto z_{1}+\sum A_{\alpha} z_{\alpha}$. Consideration of $q\left(z_{1}+\sum A_{\alpha} z_{\alpha}\right)$ shows that we must solve

$$
\begin{aligned}
& 2 \gamma \bar{A}_{\alpha}+A_{\alpha}=-b_{\alpha} \\
& \bar{A}_{\alpha}+2 \gamma A_{\alpha}=-\bar{b}_{\alpha}
\end{aligned}
$$

which is always possible if $\gamma \neq 1 / 2$. The remaining terms in $q_{1}$ and $q_{2}$ are eliminated by the change $z_{n} \mapsto z_{n}-\Sigma\left(a_{\alpha} z_{\alpha} z_{1}+c_{\alpha \beta} z_{\alpha} z_{\beta}\right)$. Next we eliminate the quadratic terms in $f_{\alpha}$, which have the form $q_{\alpha}+q_{1 \alpha}+q_{2 \alpha}$,

$$
\begin{aligned}
& q_{\alpha}=a_{a} z_{1}^{2}+b_{a} z_{1} \bar{z}_{1}+\bar{a}_{\alpha} \bar{z}_{1}^{2}, \quad b_{\alpha}=\bar{b}_{a}, \\
& q_{1 \alpha}=2 \operatorname{Re} \sum c_{\alpha \beta} x_{\beta} z_{1}, \\
& q_{2 \alpha}=\sum c_{\alpha \beta \sigma} x_{\beta} x_{\sigma}, \quad c_{\alpha \beta \sigma}=\bar{c}_{\alpha \beta \sigma} .
\end{aligned}
$$

The $z_{1} \bar{z}_{1}$-term in $q_{\alpha}$ is removed by $z_{\alpha} \mapsto z_{\alpha}+i b_{\alpha} z_{n}$. With this term zero, the remaining quadratic terms are removed by $z_{\alpha} \rightarrow z_{\alpha}+2 i \sum c_{\alpha \beta} z_{\beta} z_{1}+i \sum c_{\alpha \beta} z_{\beta} z_{\alpha}+2 i a_{a} z n$.

From this point on we assume that $M$ has the form (1.16) with

$$
\begin{align*}
& F=q\left(z_{1}, \bar{z}_{1}\right)+H\left(z_{1}, \bar{z}_{1}, x\right), \quad H=O\left(|z|^{3}\right), \\
& f_{a}=h_{a}\left(z_{1}, \bar{z}_{1}, x\right)=\bar{h}_{a}\left(\bar{z}_{1}, z_{1}, x\right), \quad h_{\alpha}=O\left(|z|^{3}\right),  \tag{1.17}\\
& q=\gamma z_{1}^{2}+z_{1} \bar{z}_{1}+\gamma \bar{z}_{1}^{2}, \quad 0 \leqslant \gamma<\infty, \gamma \neq \frac{1}{2} .
\end{align*}
$$

Before continuing let us examine the locus $N$ of those points near the origin at which $M$ has a complex tangent. We set $r^{0}=F-z_{n}, r^{\alpha}=f_{\alpha}-y_{\alpha}$, and

$$
\Delta=\frac{\partial\left(r^{0}, r^{a}, \tilde{r}^{0}\right)}{\partial\left(z_{1}, z_{\beta}, z_{n}\right)}
$$

the Jacobian determinant. Then $N$ is given by (1.16) together with $\Delta=\bar{\Delta}=0$. In view of (1.17) $\Delta= \pm(i / 2)^{n-2}\left(\bar{z}_{1}+2 \gamma z_{1}\right)+\ldots$. The condition $\gamma \neq 1 / 2$ allows us to solve $\Delta=\bar{\Delta}=0$ explicitly for $z_{1}$ and $\bar{z}_{1}: z_{1}=\psi(x), \bar{z}_{1}=\bar{\psi}(x)$. Thus if $\gamma \neq 1 / 2$, then $N$ is a totally real $(n-2)$ dimensional manifold lying on $M$.

Now we complexify the manifold $M$ by replacing $\bar{z}$ by $w$ in (1.13) to get a complex analytic $n$-dimensional submanifold $\mathfrak{M}$ in $\mathbf{C}^{2 n}$,

$$
\begin{align*}
z_{n} & =F\left(z_{1}, w_{1}, x\right), \quad 2 x_{\alpha}=z_{a}+w_{a} \\
w_{n} & =\tilde{F}\left(w_{1}, z_{1}, x\right),  \tag{1.18}\\
z_{a}-w_{a} & =2 i f_{a}\left(z_{1}, w_{1}, x\right)=2 i f_{a}\left(w_{1}, z_{1}, x\right) .
\end{align*}
$$

We note that these equations imply

$$
\begin{array}{r}
z_{\alpha}=x_{\alpha}+i f_{\alpha}\left(z_{1}, w_{1}, x\right) \\
w_{\alpha}=x_{\alpha}-i f_{\alpha}\left(z_{1}, w_{1}, x\right) . \tag{1.19}
\end{array}
$$

The variables $\left(z_{1}, w_{1}, x\right)$ will be used as complex coordinates on $\mathfrak{M}$. The two projections $\pi_{1}(z, w)=z$ and $\pi_{2}(z, w)=w$, when restricted to $\mathfrak{M}$ have the form

$$
\begin{aligned}
& \pi_{1}\left(z_{1}, w_{1}, x\right)=\left(z_{1}, x_{\alpha}+i f_{\alpha}\left(z_{1}, w_{1}, x\right), F\left(z_{1}, w_{1}, x\right)\right) \\
& \pi_{2}\left(z_{1}, w_{1}, x\right)=\left(w_{1}, x_{\alpha}-i f_{\alpha}\left(z_{1}, w_{1}, x\right), \bar{F}\left(w_{1}, z_{1}, x\right)\right)
\end{aligned}
$$

Since $\mathfrak{M}$ comes from a real submanifold $M$, the reflection $\varrho(z, w)=(\bar{w}, \bar{z})$ preserves $\mathfrak{M}$ and induces the anti-holomorphic involution $\varrho\left(z_{1}, w_{1}, x\right)=\left(\bar{w}_{1}, \bar{z}_{1}, \bar{x}\right)$.

Again the case $\gamma=0$ is exceptional, so we assume that $0<\gamma<\infty, \gamma \neq 1 / 2$. We define a holomorphic involution $\tau_{1}(z, w)=\left(z^{\prime}, w^{\prime}\right)$ on $\mathfrak{M}$ by $w=w^{\prime}$, which amounts to the equations

$$
\begin{gathered}
q\left(z_{1}^{\prime}, w_{1}\right)-q\left(z_{1}, w_{1}\right)=\bar{H}\left(w_{1}, z_{1}, x\right)-\bar{H}\left(w_{1}, z_{1}^{\prime}, x^{\prime}\right) \\
x_{\alpha}^{\prime}-i h_{\alpha}\left(z_{1}^{\prime}, w_{1}, x^{\prime}\right)=x_{\alpha}-i h_{\alpha}\left(z_{1}, w_{1}, x\right)
\end{gathered}
$$

By the implicit function theorem we get

$$
\begin{align*}
z_{1}^{\prime} & =-z_{1}-\frac{1}{\gamma} w_{1}+K+\left(z_{1}, w_{1}, x\right) \\
\tau_{1}: w_{1}^{\prime} & =w_{1}  \tag{1.20}\\
x_{a}^{\prime} & =x_{\alpha}+L_{a}\left(z_{1}, w_{1}, x\right)
\end{align*}
$$

for certain functions $K, L_{\alpha}$ of second order. The condition $\tau_{1}^{2}=$ id gives $K \circ \tau_{1}=K$, $L_{a} \circ \tau_{1}=-L_{\alpha}$. From $\tau_{2}=\varrho \tau_{1} \varrho$, we have

$$
\begin{align*}
z_{1}^{\prime} & =-z_{1} \\
\tau_{2}: w_{1}^{\prime} & =-\frac{1}{\gamma} z_{1}-w_{1}+\bar{K}\left(w_{1}, z_{1}, x\right)  \tag{1.21}\\
x_{a}^{\prime} & =x_{a}+\bar{L}_{a}\left(w_{1}, z_{1}, x\right) .
\end{align*}
$$

$\pi_{1}$ is a two-fold branched covering with covering transformation $\tau_{2}$. To find the branch locus consider the Jacobian determinant

$$
\Delta=\frac{\partial\left(z_{1}, z_{\alpha}, z_{n}\right)}{\partial\left(z_{1}, x, w_{1}\right)}
$$

where $z_{\alpha}, z_{n}$ are given by the first equations in (1.18) and (1.19). Since

$$
\frac{\partial\left(\Delta, x_{a}+i f_{a}\right)}{\partial\left(w_{1}, x_{\beta}\right)}(0)=\Delta_{w_{1}}(0)=2 \gamma
$$

$\Delta=0$ and the first equation in (1.19) can be used to eliminate $\left(w_{1}, x\right)$ in the first equation of (1.18). This shows that the branch loci of $\pi_{1}$ and $\pi_{2}=c \pi_{1} \varrho \varrho$ are smooth analytic hypersurfaces in $z$ - and $w$-space, respectively. The same argument as in the two dimensional case shows that an analytic function $f=f\left(z_{1}, w_{1}, x\right)$ is the trace on $\mathfrak{M}$ of a function holomorphic in $z$ if and only if $f \circ \tau_{2}=f$. Thus the study of the $n$-manifold $M$ also leads to consideration of a triple of involutions ( $\tau_{1}, \tau_{2}, \varrho$ ).

## 2. Quadrics and linear involutions

In this section we consider the case in which $M^{n} \subset \mathbf{C}^{n}$ is the quadric $Q_{\gamma}$

$$
\begin{align*}
Q_{\gamma}: & \begin{array}{l}
z_{n} \\
y_{\alpha}
\end{array}=q\left(z_{1}, \bar{z}_{1}\right) \equiv q_{\gamma}\left(z_{1}, \bar{z}_{1}\right)  \tag{2.1}\\
q_{\gamma} & =\gamma z_{1}^{2}+z_{1} \bar{z}_{1}+\gamma \bar{z}_{1}^{2}, \quad 0 \leqslant \gamma<\infty-1, \\
q_{\infty} & =z_{1}^{2}+\bar{z}_{1}^{2} .
\end{align*}
$$

The coordinates are as in (1.15). This will be a prelude to the study of the general manifolds of section 1 . The cases $\gamma=0,1 / 2, \infty$ are exceptional and enter the discussion only in a minor way. We also consider the complex quadrics

$$
\mathcal{Q}_{\gamma}: \begin{align*}
& z_{n}=w_{n}=q\left(z_{1}, w_{1}\right)  \tag{2.2}\\
& z_{a}=w_{a}, \quad 2 x_{\alpha}=\left(z_{a}+w_{a}\right)
\end{align*}
$$

where $q=q_{\gamma}$ is of the same form, but with $\gamma$ complex. $\Omega_{\gamma}$ may come from a $Q_{\gamma}$ by complexification.

The projections $\pi_{1}(z, w)=z, \pi_{2}(z, w)=w$ restricted to $\mathscr{Q}_{\gamma}$ are given by the quadratic mappings

$$
\begin{aligned}
& \pi_{1}\left(z_{1}, w_{1}, x\right)=\left(z_{1}, x, q\left(z_{1}, w_{1}\right)\right) \\
& \pi_{2}\left(z_{1}, w_{1}, x\right)=\left(w_{1}, x, q\left(z_{1}, w_{1}\right)\right)
\end{aligned}
$$

If $\gamma=0$, then $\pi_{1}$ collapses the lines $z_{1}=0, x=$ const. to points and is otherwise one-toone. If $\gamma \neq 0$, then $\pi_{1}$ and $\pi_{2}$ are two-fold branched coverings having as covering transformations two linear involutions $\tau_{2}$ and $\tau_{1}$. The $w$-planes cut $\mathcal{Q}_{\gamma}$ in the pointpairs of the involution $\tau_{1}$, while the $z$-planes cut $Q_{\gamma}$ in those of $\tau_{2}$. Letting $X=\left(z_{1}, w_{1}, x\right)^{t}$ be a column coordinate vector we have, as in (1.20, 1.21),

$$
\begin{equation*}
\tau_{j}(X)=T_{j} X, T_{j}^{2}=I, \quad j=1,2 \tag{2.3}
\end{equation*}
$$

where

$$
T_{1}=\left[\begin{array}{rcc}
-1 & -\gamma^{-1} & 0  \tag{2.4}\\
0 & 1 & 0 \\
0 & 0 & I
\end{array}\right], \quad T_{2}=\left[\begin{array}{crr}
1 & 0 & 0 \\
-\gamma^{-1} & -1 & 0 \\
0 & 0 & I
\end{array}\right], \quad I=I_{n-2}
$$

These formulas are also valid for $\gamma=\infty$. Each $\tau_{j}$ has a one-dimensional ( -1 )-eigenspace, so is a reflection in a hyperplane $E_{j}$. The cases $\gamma= \pm 1 / 2$ correspond to $E_{1}=E_{2}$. If $\gamma=\infty$ then a $(-1)$-eigenvector of $\tau_{1}$ is a $(+1)$-eigenvector of $\tau_{2}$, and conversely. Aside from these exceptional cases $E=E_{1} \cap E_{2}$, the space of points fixed by both $\tau_{1}$ and $\tau_{2}$, has dimension $n-2$, and $\tau_{1}$ and $\tau_{2}$ have no other common eigenvectors. The plane $F: x=0$ is invariant under both $\tau_{1}$ and $\tau_{2}$. If $\mathfrak{Q}_{\gamma}$ is the complexification of $Q_{\gamma}$, then $\mathfrak{Q}_{\gamma}$ carries the linear anti-holomorphic involution $\varrho\left(z_{1}, w_{1}, x\right)=\left(\bar{w}_{1}, \bar{z}_{1}, \bar{x}\right)$ and $\varrho \tau_{1}=\tau_{2} \varrho$. $\varrho$ preserves both $F$ and $E$, and $E=N+i N$ where $N$ is pointwise fixed by $\varrho . N$ is the locus of points at which $Q_{\gamma}$ has a complex tangent.

We now turn to the theory of a pair of holomorphic involutions on $C^{n}$, which we assume to be given in the form (2.3). First we consider the case $n=2$. The complex 2 by 2 matrices $T_{j}$ are assumed to satisfy

$$
T_{j}^{2}=I, \quad \operatorname{det} T_{j}+1=\operatorname{tr} T_{j}=0
$$

Also, we require that $T_{1}$ and $T_{2}$ have no eigenvectors in common. The mapping $\varphi=\tau_{1} \tau_{2}$ has the matrix form

$$
\begin{equation*}
\varphi(X)=\Phi X, \quad \Phi=T_{1} T_{2}, \quad \operatorname{det} \Phi= \pm 1 \tag{2.5}
\end{equation*}
$$

Lemma 2.1. Let the linear transformations $\tau_{1}, \tau_{2}, \varphi$ on $\mathbf{C}^{2}$ be as just described. Then $\varphi$ is diagonalizable with distinct eigenvalues $\mu, \mu^{-1}, \mu^{2} \neq 1$. If $\left(e_{1}, e_{2}\right)$ is a basis for which

$$
\varphi\left(e_{1}\right)=\mu e_{1}, \quad \varphi\left(e_{2}\right)=\mu^{-1} e_{2}
$$

then

$$
\tau_{j}\left(e_{1}\right)=\lambda_{j}^{-1} e_{2}, \quad \tau_{j}\left(e_{2}\right)=\lambda_{j} e_{1}
$$

where $\lambda_{1} \lambda_{2}^{-1}=\mu$. The eigenvectors $\left(e_{1}, e_{2}\right)$ may be chosen so that $\lambda_{1}=\lambda_{2}^{-1} \equiv \lambda, \lambda^{4} \neq 1$, and are determined up to $\left(e_{1}, e_{2}\right) \mapsto\left(\alpha e_{1}, \pm \alpha e_{2}\right)$ or $\left(e_{1}, e_{2}\right) \mapsto\left(e_{2}, e_{1}\right) . \operatorname{tr} \varphi=\lambda^{2}+\lambda^{-2}$ is an invariant of $\tau_{1}, \tau_{2}$.

Proof. Let $v$ be an eigenvector of $\varphi$ with eigenvalue $\mu$. Then

$$
\varphi(v)=\tau_{1} \tau_{2}(v)=\mu v, \quad \text { or } \quad \tau_{2}(v)=\mu \tau_{1}(v)=\mu \varphi \tau_{2}(v)
$$

Thus $\tau_{2}(v)$ is also an eigenvector of $\varphi$ with eigenvalue $\mu^{-1} . \tau_{2}(v)$ and $v$ are independent, since a relation $c v=\tau_{2}(v)=\mu \tau_{1}(v)$ would imply that $v$ is a common eigenvector of both $\tau_{1}$ and $\tau_{2}$. If $\mu=\mu^{-1}$, then $\varphi= \pm \mathrm{id}$, or $\tau_{2}= \pm \tau_{1}$, which again implies a common eigenvector. Relative to the basis $e_{1}=v, e_{2}=\tau_{2}(v), \tau_{1}$ and $\tau_{2}$ satisfy the above relations with $\lambda_{1}=\mu$, $\lambda_{2}=1$. The change of eigenvectors of $\varphi,\left(e_{1}, e_{2}\right) \mapsto\left(\alpha e_{1}, \beta e_{2}\right)$, results in $\lambda_{j} \mapsto \beta \lambda_{j} \alpha^{-1}$. Hence, we can arrange that $\lambda_{1}=\lambda_{2}^{-1}=\lambda, \lambda^{2}=\mu$. We must then restrict to $\alpha= \pm \beta$. Q.E.D.

Now suppose that $\tau_{1}$ and $\tau_{2}$ satisfy $\varrho \tau_{1}=\tau_{2} \varrho$, for some linear anti-holomorphic involution $\varrho, \varrho^{2}=\mathrm{id}$,

$$
\varrho(X)=P \bar{X}, \quad P \bar{P}=1, \quad P \bar{T}_{1}=T_{2} P
$$

Again let $v$ be an eigenvector of $\varphi$ with eigenvalue $\mu$. Then, since $\varphi \varrho \varphi=\tau_{1} \varrho \tau_{2}=\varrho$,

$$
\varrho(v)=\varphi \varrho(\mu v)=\bar{\mu} \varphi \varrho(v)
$$

so that $\varrho(v)$ is an eigenvector of $\varphi$ with eigenvalue $\bar{\mu}^{-1}$. Hence, either
(i) $\mu=\bar{\mu}$,
or
(ii) $\mu \bar{\mu}=1$.

Suppose $\mu$ is real and let $e_{1}, e_{2}$ be eigenvectors of $\varphi$ as in Lemma 2.1. Since $\bar{\mu}^{-1}$ is the eigenvalue $\mu^{-1}$ and $\varrho^{2}=$ id, we have

$$
\varrho\left(e_{1}\right)=a e_{2}, \quad \varrho\left(e_{2}\right)=\bar{a}^{-1} e_{1}
$$

From $\varrho \tau_{1}\left(e_{1}\right)=\tau_{2} \varrho\left(e_{1}\right)$ we get

$$
\lambda_{1} \bar{\lambda}_{2} a \bar{a}=1
$$

It follows that

$$
\begin{equation*}
\mu=\lambda_{1} \lambda_{2}^{-1}=a \bar{a} \lambda_{1} \lambda_{1}>0 \tag{2.6}
\end{equation*}
$$

The change $\left(e_{1}, e_{2}\right) \mapsto\left(\alpha e_{1}, \beta e_{2}\right)$ results in $\left(a, \lambda_{j}\right) \mapsto\left(\bar{\alpha} a \beta^{-1}, \beta \lambda_{j} \alpha^{-1}\right)$. To make $a=1$, $\lambda_{1} \lambda_{2}=1$ by such a change, we require

$$
\beta=a \bar{\alpha}, \quad \alpha^{2} \beta^{-2}=\lambda_{1} \lambda_{2}
$$

The second condition is $(\alpha / \bar{\alpha})^{2}=a^{2} \lambda_{1} \lambda_{2}=a \lambda_{2}\left(\bar{a} \bar{\lambda}_{2}\right)^{-1}$. Since this last term has modulus one, such an $\alpha$ exists. with this normalization, $\lambda_{1}=\lambda_{2}^{-1}=\lambda=\bar{\lambda}$. We must now restrict to
$\beta=\bar{\alpha}, \alpha^{2}=\beta^{2}$. If we arrange that $\lambda>0$, then we must have $\beta=\alpha=\bar{\alpha}$. By interchanging $e_{1}$ and $e_{2}$ we may take $\lambda>1$.

Now suppose $\mu \bar{\mu}=1$, so that $\varphi$ has eigenvalues $\mu=\tilde{\mu}^{-1}$ and $\mu^{-1}$. Then $\varrho\left(e_{1}\right)=a e_{1}$, and similarly, $\varrho\left(e_{2}\right)=b e_{2}$. From $\varrho^{2}=$ id we get $a \bar{a}=b \bar{b}=1$. The above change of eigenvectors results in $(a, b) \mapsto\left(\bar{\alpha} a \alpha^{-1}, \bar{\beta} b \beta^{-1}\right)$. Therefore we can choose $e_{1}, e_{2}$ so that $a=b=1$. Now we must restrict to $\alpha$ and $\beta$ real. The relation $\varrho \tau_{1}\left(e_{1}\right)=\tau_{2} \varrho\left(e_{1}\right)$ gives $\bar{\lambda}_{1}=\lambda_{2}$. Hence by choice of real $\alpha$ and $\beta$, we can make $\lambda_{1} \lambda_{2}=1$. Thus, $\lambda_{1}=\lambda_{2}^{-1}=\lambda, \lambda \bar{\lambda}=1$. We can arrange that $\operatorname{Re} \lambda>0$, then we must restrict to $\beta=\alpha=\bar{\alpha}$. By interchanging $\dot{e}_{1}$ and $e_{2}$ we can make $0<\arg \lambda<\pi / 2$.

We introduce coordinates $(\xi, \eta)$ by $X=\xi e_{1}+\eta e_{2}$, where $\left(e_{1}, e_{2}\right)$ are the above chosen eigenvectors of $\varphi$. We have proved the following.

LEMMA 2.2. Let $\tau_{1}, \tau_{2}, \varphi$ be as in Lemma 2.1 and suppose that $\varrho \tau_{1}=\tau_{2} \varrho$ for some linear anti-holomorphic involution $\varrho$. Then there exist linear coordinates $(\xi, \eta)$ in which

$$
\begin{equation*}
\tau_{1}(\xi, \eta)=\left(\lambda \eta, \lambda^{-1} \xi\right), \quad \tau_{2}(\xi, \eta)=\left(\lambda^{-1} \eta, \lambda \xi\right) \tag{2.7}
\end{equation*}
$$

Also, either
(i) $\varrho(\xi, \eta)=(\bar{\eta}, \xi)$ and $\lambda=\bar{\lambda}>1$, or
(ii) $\varrho(\xi, \eta)=(\bar{\xi}, \bar{\eta})$ and $\lambda \bar{\lambda}=1,0<\arg \lambda<\pi / 2$.

Such coordinates are determined up to $(\xi, \eta) \mapsto(\alpha \xi, \alpha \eta), \alpha=\bar{\alpha}$.
Next we consider two linear holomorphic involutions $\tau_{1}, \tau_{2}$ on $C^{n}$. We assume that each $\tau_{j}$ is a reflection in a hyperplane $E_{j}$ and that $E_{1} \neq E_{2}$. Let $E=E_{1} \cap E_{2}$ and $v_{j}$ be a (-1)-eigenvector of $T_{j}, j=1,2$. We also assume that $E, v_{1}, v_{2} \operatorname{span} \mathbf{C}^{n}$.

Lemma 2.3. (a) Let $\tau_{1}, \tau_{2}$ be involutions on $\mathrm{C}^{n}$ as just described. There exist complex linear coordinates $\xi, \eta, \zeta=\left(\zeta_{3}, \ldots, \zeta_{n}\right)$ in which

$$
\begin{equation*}
\tau_{j}(\xi, \eta, \zeta)=\left(\lambda_{j} \eta, \lambda_{j}^{-1} \xi, \zeta\right), \quad j=1,2 \tag{2.9}
\end{equation*}
$$

They may be so chosen that $\lambda_{1}=\lambda_{2}^{-1}=\lambda$ and are then determined up to replacement by ( $\alpha \xi, \pm \alpha \eta, B \zeta$ ), $\alpha \in \mathbf{C}, B \in G L(n-2, C)$, or by $(\eta, \xi, \zeta)$.
(b) If also $\varrho \tau_{1}=\tau_{2} \varrho$ for a linear anti-holomorphic involution $\varrho$, then these coordinates can be further specialized so that either
(i) $\varrho(\xi, \eta, \xi)=(\bar{\eta}, \xi, \bar{\xi})$ and $\lambda>1$, or
(ii) $\varrho(\xi, \eta, \zeta)=(\bar{\xi}, \bar{\eta}, \bar{\zeta})$ and $\lambda \bar{\lambda}=1,0<\arg \lambda<\pi / 2$.

They are then determined up to replacement by $(\alpha \xi, \alpha \eta, B \zeta), \alpha \in \mathbf{R}, B \in G L(n-2, \mathbf{R})$.
Proof. (a) Let $F$ be the space spanned by $v_{1}$ and $v_{2}$ so that $\mathbf{C}^{n}=F \oplus E$. We claim that $F$ is invariant under both $\tau_{1}$ and $\tau_{2}$. Since $\tau_{1} v_{1}=-v_{1}$, consider

$$
\tau_{1} v_{2}=\alpha v_{1}+\beta v_{2}+w, \quad w \in E
$$

We must show $\boldsymbol{w}=\mathbf{0}$. Since $\tau_{1}^{2}=\mathrm{id}$,

$$
v_{2}=-\alpha v_{1}+\beta \tau_{1} v_{2}+w=(\beta-1) \alpha v_{1}+\beta^{2} v_{2}+(\beta+1) w
$$

If $w \neq 0$, then $\left(v_{1}, v_{2}, w\right)$ are independent so $\beta=-1$, and $\alpha=0$. This implies that

$$
\tau_{1}\left(v_{2}-\frac{1}{2} w\right)=-\left(v_{2}-\frac{1}{2} w\right)
$$

hence $v_{2}-w / 2=c v_{1}$, which contradicts independence. Hence, $w=0$, and $\tau_{1}(F)=F$. A similar argument shows that $\tau_{2}(F)=F$. Let $\tau_{j}^{\prime}$ be the restriction of $\tau_{j}$ to $F$. Then $\operatorname{det} \tau_{j}^{\prime}=-1$, and by the condition on $v_{1}, v_{2}$ and $E, \tau_{1}^{\prime}$ and $\tau_{2}^{\prime}$ can have no common eigenvector. Hence, we may apply Lemma 2.1 to $\tau_{1}^{\prime}, \tau_{2}^{\prime}$, to get basis vectors $e_{1}, e_{2}$ of $F$. We let $e_{3}, \ldots, e_{n}$ be any basis of $E$, and $(\xi, \eta, \xi)$ coordinates relative to $e_{1}, \ldots, e_{n}$.
(b) We first show that $\varrho$ leaves $E$ invariant. If $\tau_{j} w=w$, then $\tau_{j} \varrho(w)=\varrho(w)$ follows from $\tau_{1} \varrho=\varrho \tau_{2}$, hence $\varrho(E)=E$. Let $N$ be the totally real fixed point set of $\varrho$ on $E, E=N+i N$. Choose the coordinates $\zeta$ on $E$ so that $\varrho: \zeta \mapsto \xi$. We next show that $F$ is invariant under $\varrho$. To see this note that

$$
\begin{aligned}
& \tau_{1} \varrho(F)=\varrho \tau_{2}(F)=\varrho(F) \\
& \tau_{2} \varrho(F)=\varrho \tau_{1}(F)=\varrho(F)
\end{aligned}
$$

Hence, $\varrho(F)$ is invariant under both $\tau_{1}$ and $\tau_{2}$. Relative to a basis compatible with the decomposition $\varrho(F) \oplus E=\varrho(F) \oplus \varrho(E)=\mathbf{C}^{n}$ it is easy to see that $\operatorname{det} \tau_{j}^{\prime}=-1$, where $\tau_{j}^{\prime}=\left.\tau_{j}\right|_{\varrho(F)}$. So $\tau_{j}^{\prime}$ has a $(-1)$-eigenvector $u_{j}$ in $\varrho(F)$. By the assumption made on $\tau_{1}, \tau_{2}$, we must have $u_{j}=c v_{j}$. It follows that $u_{1}, u_{2}$ are independent and $\varrho(F)=F$. We now apply Lemma 2.2 to $\tau_{1}^{\prime}, \tau_{2}^{\prime}$.
Q.E.D.

Given involutions $\tau_{1}, \tau_{2}, \varrho$ as in $(2.7,2.8)$ in canonical coordinates $(\xi, \eta, \zeta)$ $\left(\lambda_{1}=\lambda_{2}^{-1}=\lambda\right.$ ), we shall construct a quadric $Q_{\gamma}$. For this we must construct the "holomorphic" coordinates $z$ and the "anti-holomorphic" coordinates $w$. Linear combina-
tions of the $\zeta_{\alpha}$ are invariant under both $\tau_{1}$ and $\tau_{2}$. Aside from these the most general linear functions invariant under $\tau_{2}$ and $\tau_{1}$, respectively, are

$$
\begin{aligned}
z_{1} & =b(\lambda \xi+\eta) \\
w_{1} & =a(\xi+\lambda \eta)
\end{aligned}
$$

where $a$ and $b$ are complex constants. We should also choose $z_{1}$ and $w_{1}$ so that $\varrho$ corresponds to $\left(z_{1}, w_{1}, x\right) \mapsto\left(\bar{w}_{1}, \bar{z}_{1}, \bar{x}\right)$. If $\lambda=\bar{\lambda}$ we need $a=\bar{b}$, while if $\lambda \bar{\lambda}=1$ we need $a \lambda=\bar{b}$. Thus for the two cases in (2.10) we take
(i) $\begin{aligned} z_{1} & =b(\lambda \xi+\eta) \\ w_{1} & =\bar{b}(\xi+\lambda \eta)\end{aligned}$,
(ii) $\begin{gathered}z_{1}=b(\lambda \xi+\eta) \\ w_{1}=\bar{d} \bar{\lambda}(\xi+\lambda \eta) .\end{gathered}$

The quadratic functions invariant under both $\tau_{1}$ and $\tau_{2}$ are linear combinations of $\xi \eta$ and $\zeta_{\alpha} \zeta_{\beta}$.

We want to choose $b$ so that $q\left(z_{1}, w_{1}\right)$ is a multiple of $\xi \eta$, for some $q$ of the form (2.1). In case (ii) this requires that

$$
b^{2} \lambda^{2}+b^{2} \lambda^{2}=b^{2}+\bar{b}^{2}
$$

Taking $b \bar{b}=1$, we get $b^{4}=\lambda^{-2}$.
Hence, in both cases, we arrive at

$$
\begin{align*}
& z_{1}=i \lambda^{-1 / 2}(\lambda \xi+\eta), \\
& w_{1}=-i \lambda^{-1 / 2}(\xi+\lambda \eta),  \tag{2.11}\\
& z_{a}=w_{\alpha}=x_{\alpha}=\zeta_{\alpha} .
\end{align*}
$$

It follows that

$$
\begin{align*}
& q \equiv z_{1} w_{1}+\gamma\left(z_{1}^{2}+w_{1}^{2}\right)=\gamma^{-1}\left(1-4 \gamma^{2}\right) \xi \eta \\
& \gamma=\left(\lambda+\lambda^{-1}\right)^{-1}>0 \tag{2.12}
\end{align*}
$$

We define

$$
\begin{equation*}
z_{n}=w_{n}=\gamma^{-1}\left(1-4 \gamma^{2}\right) \xi \eta . \tag{2.13}
\end{equation*}
$$

$\tau_{1}, \tau_{2}, \varrho$ are the involutions induced on the surface (2.12, 2.13). We may write the relation between $\gamma$ and $\lambda$ as

$$
\begin{equation*}
\gamma \lambda^{2}-\lambda+\gamma=0 . \tag{2.14}
\end{equation*}
$$

If $\lambda \neq \pm 1$ is real, (2.14) has two distinct real roots, $\lambda, \lambda^{-1}$. It follows that $\gamma<1 / 2$ and the surface is elliptic. If $\lambda \bar{\lambda}=1$ and $0<\operatorname{Re} \lambda<1$, then $\gamma=(1 / 2)(\operatorname{Re} \lambda)^{-1}>1 / 2$ and the surface is hyperbolic.

Conversely, given $\tau_{1}$ and $\tau_{2}$ with the matrices (2.4), we can use (2.11) to define the canonical coordinates $(\xi, \eta, \zeta)$. We only have to find $\lambda . \lambda^{2}=\mu$ is an eigenvalue of $\varphi=\tau_{1} \tau_{2}=\tau_{1} \varrho \tau_{1} \varrho$. In terms of matrices $\Phi=T_{1} T_{2}=T_{1} P \bar{T}_{1} \bar{P}=\left(-T_{1} P\right)^{2}$, where $T_{i}$ are given by (2.4) and

$$
P=\left[\begin{array}{lll}
0 & 1 & 0 \\
1 & 0 & 0 \\
0 & 0 & 1
\end{array}\right] .
$$

The eigenvalues of $-T_{1} P$ are given by (2.14) together with $\lambda=-1$.
The mapping (2.11), (2.13) has an interesting geometric interpretation. In the elliptic case the relation $w_{1}=\bar{z}_{1}$ corresponds to $\eta=\xi$. Under (2.11) the ellipses $q\left(z_{1}, z_{1}\right)=c>0$ are mapped to the circles $\xi \bar{\xi}=c^{\prime}>0$. In fact ( $2.11,2.13$ ) maps $Q_{\gamma}$ to $Q_{0}$. In the hyperbolic case the relation $w_{1}=\bar{z}_{1}$ corresponds to $\xi=\xi, \eta=\bar{\eta}$. The hyperbolas $q\left(z_{1}, \bar{z}_{1}\right)=c$ are mapped to the standard hyperbolas $\xi \eta=c^{\prime} . Q_{\gamma}$ is mapped to $Q_{\infty}$ by (2.11, 2.13). Of course, $(2.11,2.13)$ is not holomorphic in the usual sense.

The linear map $\varphi=\tau_{1} \tau_{2}$,

$$
\varphi(\xi, \eta, \zeta)=\left(\mu \xi, \mu^{-1} \eta, \zeta\right)
$$

leaves fixed the linear space $\xi=\eta=0$, i.e. $z_{1}=w_{1}=0$. As mentioned above this is the complexification of the space of those points on $Q_{\gamma}$ having complex tangents. When $n=2, \varphi$ has an isolated fixed point which is hyperbolic ( $\mu>0$ ) if $Q_{\gamma}$ is elliptic, and elliptic $(\mu \bar{\mu}=1)$ when $Q_{\gamma}$ is hyperbolic. $\varphi$ may be interpolated by the flow

$$
\begin{equation*}
\varphi^{t}(\xi, \eta, \zeta)=\left(e^{t v} \xi, e^{-t v} \eta, \zeta\right), \tag{2.15}
\end{equation*}
$$

where $e^{\nu}=\mu, \varphi^{1}=\varphi$, and either
(i) $\nu=\bar{v}$
or
(ii) $v+\bar{v}=0$.
$\varphi^{t}$ preserves the family of complex conics $\xi \eta=$ const., i.e. $q\left(z_{1}, w_{1}\right)=c$, on $\mathfrak{\Re}_{\gamma}$. If $c$ is real the complex conic meets $Q_{\gamma}$ in a real conic, which may be degenerate. $\varphi$ does not preserve $Q_{\gamma}$ since $\varrho^{-1} \varphi \varrho=\varphi^{-1} \neq \varphi$, if $\gamma \neq \infty$. If we allow complex $t$, then $\varphi^{\prime} \varrho=\varrho \varphi^{-i}$ in both cases. Thus $\varphi^{t}$ commutes with $\varrho$ precisely when $t+\bar{t}=0$. The orbits of $\varphi^{t}$ on $Q_{\gamma}$ for $t+\bar{t}=0$ are the real conics. The infinitesimal generator is a vector field on $Q_{\gamma}$ tangent to these curves.

As mentioned in section $1, \varphi^{2}=\left[\tau_{1}, \tau_{2}\right]=$ id is a direct analogue of the vanishing of the Levi-form on a real hypersurface. Among the quadrics $Q_{\gamma}$ this happens only when $\gamma=\infty . Q_{\infty}$ is the intersection of $t w o$ Levi-flat hypersurfaces $\operatorname{Re}\left(z_{2}-2 z_{1}^{2}\right)=\operatorname{Im} z_{2}=0$. A weaker condition is that $\varphi$ should be nilpotent. This happens precisely when $\lambda$ is a root of unity and causes difficulties for the normal form in section 3 . The eigenvalues of $\varphi$ are multiple precisely when $\mu= \pm 1$. It is an interesting fact that $\varphi$ is diagonalizable ( $\varphi=-I$ ) for $\mu=-1$, i.e. $\gamma=\infty$, while for $\mu=+1$, i.e. $\gamma=1 / 2, \varphi$ is not diagonalizable.

Finally, we make a remark on the automorphism group of $Q_{\gamma}, \gamma \neq 0,1 / 2, \infty$. It is clear that the holomorphic map

$$
\begin{align*}
& z_{1}^{\prime}=a\left(z_{n}, z_{\alpha}\right) z_{1}, \quad a=\bar{a}, b_{a}=\bar{b}_{\alpha}, \\
& z_{n}^{\prime}=a^{2}\left(z_{n}, z_{\alpha}\right) z_{n}, \quad a \neq 0, \operatorname{det} b^{\prime} \neq 0,  \tag{2.16}\\
& z_{\alpha}^{\prime}=b_{a}\left(z_{n}, z_{\beta}\right), \quad \text { for } z=0,
\end{align*}
$$

preserves $Q_{\gamma}$. Via the mapping (2.11) this corresponds to

$$
(\xi, \eta, \zeta) \mapsto\left(a(\xi \eta, \zeta) \xi, a(\xi \eta, \zeta) \eta, b_{a}(\xi \eta, \zeta) \zeta_{\alpha}\right),
$$

which is an automorphism of the set of involutions $\tau_{1}, \tau_{2}, \varrho$. In the next section we shall use this to show the most general self transformation of $Q_{\gamma}$ is of the form (2.16) where $a$ and $b_{\alpha}$ are arbitrary real formal power series, if $\gamma$ is not exceptional.

## 3. The formal theory of a pair of involutions

The considerations of section 1 have led us to a pair of holomorphic involutions $\tau_{1}, \tau_{2}$ defined in a neighborhood of a fixed point on a complex manifold $\mathfrak{M}$. In this and the following section we assume $\mathfrak{M}=\mathbf{C}^{n}$, with coordinates $x, y, z=\left(z_{a}\right), 2 \leqslant \alpha \leqslant n-1$, and that the origin is the fixed point. We now ask for a new coordinate system $\xi, \eta, \zeta=\left(\zeta_{a}\right)$ in which $\tau_{1}, \tau_{2}$ take a particularly simple form, a so-called normal form. We shall first discuss the normal form in the realm of formal power series on a purely algebraic level. Later, in the next section, we discuss the question of convergence. The case in which $\tau_{1}$ and $\tau_{2}$ are intertwined by an anti-holomorphic involution $\varrho$ will also be considered.

One may proceed directly with the mappings $\tau_{j}$; however, we shall base our analysis on the mapping $\varphi=\tau_{1} \tau_{2}$. As in the linear case (section 2 ) we shall normalize $\varphi$ and then show that this forces a normalization of $\tau_{1}$ and $\tau_{2}$.

By the results of section 2 we may take $\tau_{j}$ of the form

$$
\begin{align*}
x^{\prime} & =\lambda_{j} y+p_{j}(x, y, z) \\
\tau_{j}: y^{\prime} & =\lambda_{j}^{-1} x+q_{j}(x, y, z), \quad j=1,2 .  \tag{3.1}\\
z_{\alpha}^{\prime} & =z_{\alpha}+r_{j \alpha}(x, y, z)
\end{align*}
$$

Then $\varphi$ has the form

$$
\begin{align*}
x^{\prime} & =\mu x+f(x, y, z) \\
\varphi: y^{\prime} & =\mu^{-1} y+g(x, y, z)  \tag{3.2}\\
z_{\alpha}^{\prime} & =z_{\alpha}+h_{\alpha}(x, y, z)
\end{align*}
$$

where $\mu=\lambda_{1} \lambda_{2}^{-1}$. Here $p_{j}, q_{j}, r_{j}, f, g, h$, are formal power series vanishing to second order at the origin. We subject these mappings to the group ${ }^{(5)}{ }^{1}$ of formal transformations which agree with the identity to second order. Such a $\psi \in(s)^{1}$ has the form

$$
\begin{align*}
x & =U(\xi, \eta, \zeta)=\xi+u(\xi, \eta, \zeta) \\
\psi: y & =V(\xi, \eta, \zeta)=\eta+v(\xi, \eta, \zeta)  \tag{3.3}\\
z & =W(\xi, \eta, \zeta)=\zeta+w(\xi, \eta, \zeta)
\end{align*}
$$

where $z, w, W$ are ( $n-2$ )-vector valued and $u, v, w$ begin with quadratic terms. We call $\psi$ normalized if the power series $u, v, w$ do not contain terms of the form $\xi^{j+1} \eta^{j}, \xi^{j} \eta^{j+1}$, or $\xi^{j} \eta^{j}$, respectively, for any $j \in \mathbf{Z}^{+}$. Any formal power series $p=p(\xi, \eta, \zeta)$ may be decomposed as

$$
p=\sum_{s=-\infty}^{\infty} p_{s}(\xi, \eta, \zeta), \quad p_{s}(\xi, \eta, \zeta)=\sum_{i-j=s} \sum_{|K|=0}^{\infty} p_{i j K} \xi^{i} \eta^{j} \zeta^{K}
$$

We shall say that $p_{s}$ has type s. The normalizing conditions on $\psi$ may be expressed as

$$
\begin{equation*}
u_{1}=0, \quad v_{-1}=0, \quad w_{0}=0 \tag{3.4}
\end{equation*}
$$



$$
\psi=\psi_{0} \delta
$$

where $\psi_{0}$ is normalized and $\delta$ has the form

$$
\delta:(\xi, \eta, \zeta) \mapsto(\alpha \xi, \beta \eta, \zeta+\gamma)
$$

where $a, \beta, \gamma$ are power series in $\zeta$ and the product $\xi \eta$. If $\psi$ is convergent, so are $\psi_{0}$ and $\delta$.

Proof. We may define such $\alpha, \beta, \gamma$ by

$$
\begin{equation*}
U_{1}=\alpha(\xi \eta, \zeta) \xi, \quad V_{-1}=\beta(\xi \eta, \zeta) \eta, \quad W_{0}=\zeta+\gamma(\xi \eta, \zeta), \tag{3.5}
\end{equation*}
$$

and form $\delta$ as in the statement of the lemma. Since the transformation $p \mapsto p(\alpha \xi, \beta \eta, \zeta+\gamma)$ commutes with the projections $p \mapsto p_{s}$ it follows that $\psi_{0}=\psi \circ \delta^{-1}$ is normalized. Conversely, any such decomposition $\psi_{0}, \delta$ forces (3.5) so $\delta$ and $\psi_{0}$ are unique. It is also clear that if $\psi$ converges, so does $\delta$ and hence also $\psi_{0}$. Q.E.D.

Lemma 3.2. Let $\tau_{j}, j=1,2$ be two formal involutions given by (3.1) with $\mu=\lambda_{1} \lambda_{2}^{-1}$ not a root of unity. Then there exists a unique normalized transformation $\psi$ as in (3.3) such that relative to the coordinates $(\xi, \eta, \zeta)$

$$
\psi^{-1} \circ \tau_{j} \circ \psi:\left\{\begin{array}{l}
\xi^{\prime}=\Lambda_{j} \eta  \tag{3.6}\\
\eta^{\prime}=\Lambda_{j}^{-1} \xi, j=1,2 ; \quad \psi^{-1} \varphi \psi:\left\{\begin{array}{l}
\xi^{\prime}=M \xi \\
\xi^{\prime}=\zeta
\end{array},\left\{\begin{array}{l} 
\\
\eta^{\prime}=M^{-1} \eta \\
\xi^{\prime}=\zeta
\end{array},\right.\right.
\end{array}\right.
$$

where $M=\Lambda_{1} \Lambda_{2}^{-1}$ and the $\Lambda_{j}=\lambda_{j}+\ldots$ are formal power series in $\zeta$ and the product $\xi \eta$.
Proof. We proceed by induction on the homogeneous degree in all variables of the terms in $\psi^{-1} \tau_{j} \psi$. We assume that $\tau_{j}$ has been transformed so as to have the form (3.6) modulo terms of order $m$ and higher by a unique choice of the terms in $\psi$ of order less than $m$. It will suffice to show that the term of order $m$ in $\psi$ can be chosen uniquely so that $\psi^{-1} \tau_{j} \psi$ has the form (3.6) modulo terms of order $m+1$. Thus assume $\tau_{j}$ has the form

$$
\tau_{j}:\left\{\begin{array}{l}
x^{\prime}=\Lambda_{j} y+p_{j}+\ldots  \tag{3.7}\\
y^{\prime}=\Lambda_{j}^{-1} x+q_{j}+\ldots, \quad j=1,2 \\
z^{\prime}=z+r_{j}+\ldots
\end{array}\right.
$$

where $\Lambda_{j}=\Lambda_{j}(x y, z)$ are polynomials of degree $<m-1, p_{j}, q_{j}, r_{j}$ are homogeneous polynomials of degree $m \geqslant 2$, and the dots indicate higher order terms. Using $\tau_{j}^{2}=\mathrm{id}$ and noting $\Lambda_{j} \tau_{j}=\Lambda_{j}+O(m)$, we get

$$
\begin{align*}
& \lambda_{j} q_{j}(x, y, z)+p_{j}\left(\lambda_{j} y, \lambda_{j}^{-1} x, z\right)=0, \quad j=1,2, \\
& r_{j}(x, y, z)+r_{j}\left(\lambda_{j} y, \lambda_{j}^{-1} x, z\right)=0, \quad j=1,2 . \tag{3.8}
\end{align*}
$$

It follows that $\varphi$ has the form

$$
\varphi:\left\{\begin{array}{l}
x^{\prime}=M x+a+\ldots  \tag{3.9}\\
y^{\prime}=M^{-1} y+b+\ldots \\
z^{\prime}=z+c+\ldots
\end{array}\right.
$$

where $M=\Lambda_{1} \Lambda_{2}^{-1}$ and $a, b, c$ are the homogeneous polynomials of degree $m$ given by

$$
\begin{align*}
& a(x, y, z)=\lambda_{1} q_{2}(x, y, z)+p_{1}\left(\lambda_{2} y, \lambda_{2}^{-1} x, z\right) \\
& b(x, y, z)=\lambda_{1}^{-1} p_{2}(x, y, z)+q_{1}\left(\lambda_{2} y, \lambda_{2}^{-1} x, z\right)  \tag{3.10}\\
& c(x, y, z)=r_{2}(x, y, z)+r_{1}\left(\lambda_{2} y, \lambda_{2}^{-1} x, z\right)
\end{align*}
$$

Now let $\psi$ have the form $(3.3,3.4)$ in which $u, v, w$ are homogeneous polynomials of degree $m$. We shall choose $u, v, w$ so that $\tilde{\varphi}=\psi^{-1} \varphi \psi$ has the form given in (3.6) modulo terms of order $m+1$, and then show that automatically the $\psi^{-1} \tau_{j} \psi$ also have the form in (3.6) to the same order. Let $\tilde{\varphi}$ be as in (3.9) with $\tilde{M}=M$ and $\tilde{a}, \tilde{b}, \tilde{c}$ homogeneous of degree $m$ in $(\xi, \eta, \zeta)$. Since $M(x y, z)=M(\xi \eta, \zeta)+O(m)$, comparison of terms of degree $m$ in $\psi \tilde{\varphi}=\varphi \psi$ gives

$$
\begin{align*}
& u\left(\mu \xi, \mu^{-1} \eta, \zeta\right)-\mu u(\xi, \eta, \zeta)=(a-\tilde{a})(\xi, \eta, \zeta) \\
& v\left(\mu \xi, \mu^{-1} \eta, \zeta\right)-\mu^{-1} v(\xi, \eta, \zeta)=(b-\tilde{-})(\xi, \eta, \zeta)  \tag{3.11}\\
& w\left(\mu \xi, \mu^{-1} \eta, \zeta\right)-w(\xi, \eta, \zeta)=(c-\tilde{c})(\xi, \eta, \zeta) .
\end{align*}
$$

We wish to make $\tilde{a}_{s}=0$, for $s \neq 1, b_{s}=0$, for $s \neq-1$, and $\tilde{c}_{s}=0$, for $s \neq 0$, where $s$ indicates the type. This leads to the equations

$$
\begin{aligned}
& \left(\mu^{s}-\mu\right) u_{s}=a_{s}, \quad s \neq 1 \\
& \left(\mu^{s}-\mu^{-1}\right) v_{s}=b_{s}, \quad s \neq-1 \\
& \left(\mu^{s}-1\right) w_{s}=c_{s}, \quad s \neq 0,
\end{aligned}
$$

which clearly can be solved for $u_{s}, v_{s}, w_{s}$ since by our assumption no power of $\mu$ is unity.

For the exceptions just made the left hand sides vanish, forcing

$$
\tilde{a}_{1}=a_{1}, \quad b_{-1}=b_{-1}, \quad \tilde{c}_{0}=c_{0}
$$

The normalization (3.4) makes the solution unique. Hence, we can achieve that

$$
\begin{equation*}
a=A(x y, z) x, \quad b=B(x y, z) y, \quad c=C(x y, z) \tag{3.12}
\end{equation*}
$$

by a unique choice of the terms of order $m$ in $\psi$, if $\psi$ is normalized.

We next show this actually implies $c=0$ and $r_{j}=0, j=1,2$. To see this let $\tau_{1}, \tau_{2}, \varphi$ temporarily stand for the linear parts of these mappings. By (3.12) we have $C \tau_{1}=C \tau_{2}=C$. The third equation in (3.10) is equivalent to

$$
r_{2}+r_{1} \tau_{2}=C, \quad r_{1}+r_{2} \tau_{2}=C
$$

(The second relation follows from an application of $\tau_{2}$.) The second equation of (3.8) gives respectively for $j=1,2$,

$$
r_{2} \tau_{2}=-r_{2}, \quad r_{1} \varphi=-r_{1} \tau_{2}
$$

Hence,

$$
r_{2}-r_{1} \varphi=r_{1}-r_{2}=C
$$

or

$$
r_{1}-r_{1} \varphi=2 C
$$

Since $\mu$ is not a root of unity, this last relation implies that the terms of type $s \neq 0$ in $r_{1}$ vanish. Therefore $r_{1} \tau_{1}=r_{1}$. By (3.8) $r_{1}=0$. It follows that $r_{2}=-C$ is of type 0 , so must also vanish. We next want to show that

$$
\begin{equation*}
p_{j}(x, y, z)=P_{j}(x y, z) y, \quad q_{j}(x, y, z)=Q_{j}(x y, z) x \tag{3.13}
\end{equation*}
$$

To see this we write the first two equation of (3.10), taking into account (3.12), in the form

$$
\lambda_{1} q_{2}+p_{1} \tau_{2}=x A, \quad \lambda_{1}^{-1} p_{2}+q_{1} \tau_{2}=y B
$$

We also write the first equation in (3.8) as

$$
q_{1}=-\lambda_{1}^{-1} p_{1} \tau_{1}, \quad q_{2}=-\lambda_{2}^{-1} p_{2} \tau_{2}
$$

Eliminating $q_{1}$ and $q_{2}$ we get

$$
p_{1}-\mu p_{2}=\lambda_{2} y A, \quad p_{2}-p_{1} \varphi=\lambda_{1} y B
$$

where we have used $(x A) \tau_{2}=\lambda_{2} y A$. Eliminating first $p_{2}$ and then $p_{1}$ and using $(x A) \circ \varphi=\mu x A$, we get

$$
\begin{aligned}
& p_{1}-\mu p_{1} \varphi=y\left(\lambda_{2} A+\mu \lambda_{1} B\right) \\
& p_{2}-\mu p_{2} \varphi=y\left(\lambda_{2} \mu^{-1} A+\lambda_{1} B\right)
\end{aligned}
$$

Again, since $\mu$ is not a root of unity, these equations imply that $p_{1}$ and $p_{2}$ are of type $s=-1$ in $(x, y)$. Thus the first equation in (3.13) holds. Eliminating $p_{1}$ and $p_{2}$ via $p_{j}=-\lambda_{j} q_{j} \tau_{j}$ gives

$$
\begin{aligned}
& q_{1}-\mu^{-1} q_{1} \varphi=x\left(\mu^{-1} \lambda_{1}^{-1} A+\lambda_{2}^{-1} B\right) \\
& q_{2}-\mu^{-1} q_{2} \varphi=x\left(\lambda_{1}^{-1} A+\lambda_{2}^{-1} \mu B\right)
\end{aligned}
$$

These imply that $q_{1}$ and $q_{2}$ are of type $s=+1$. Hence, (3.13) is proved.
Returning $\tau_{j}, \varphi$ to their original meanings, we may write (3.7) as

$$
\begin{aligned}
x^{\prime} & =\left(\Lambda_{j}+P_{j}\right) y+\ldots \\
\tau_{j}: y^{\prime} & =\left(\Lambda_{j}^{-1}+Q_{j}\right) x+\ldots \\
z^{\prime} & =z+\ldots,
\end{aligned}
$$

where the dots indicate terms of order $m+1$ and higher. The relations (3.8) and (3.13) give

$$
\lambda_{j}^{-1} P_{j}+\lambda_{j} Q_{j}=0, \quad j=1,2
$$

which imply that

$$
\left(\Lambda_{j}+P_{j}\right)\left(\Lambda_{j}^{-1}+Q_{j}\right)=1+O(m)
$$

If we replace $\Lambda_{j}$ by $\Lambda_{j}+P_{j}$, then we have achieved (3.7) with the degree $m$ replaced by $m+1$. By induction we can achieve the form (3.6) for $\psi^{-1} \tau_{j} \psi$ with a unique normalized $\psi$. The form of $\psi^{-1} \varphi \psi$ follows, and the lemma is proved.
Q.E.D.

In view of the applications we wish to make to surfaces we consider the case in which $\tau_{1}$ and $\tau_{2}$ are intertwined by one of the linear anti-holomorphic involutions $\varrho$
(i) $\varrho(x, y, z)=(\bar{y}, \bar{x}, \bar{z})$,
(ii) $\varrho(x, y, z)=(\bar{x}, \bar{y}, \bar{z})$.

We have the following lemma.
Lemma 3.3. Suppose that the $\tau_{1}, \tau_{2}$ of Lemma 3.2 also satisfy $\varrho \tau_{1}=\tau_{2} \varrho$, where $\varrho$ is one of the anti-holomorphic involutions (3.14). Then the transformation $\psi$ satisfies $\psi \varrho=\varrho \psi$, and the factors $\Lambda_{1}, \Lambda_{2}$ are related by
(i) $\bar{\Lambda}_{1}(\xi \eta, \zeta)=\Lambda_{2}(\xi \eta, \zeta)^{-1}$,
(ii) $\Lambda_{1}(\xi \eta, \zeta)=\bar{\Lambda}_{2}(\xi \eta, \zeta)$.

Proof. Let $\tau_{j}^{*}$ denote the normal forms (3.6) so that $\psi \tau_{j}^{*}=\tau_{j} \psi, j=1,2$. We have $(\varrho \psi \varrho)\left(\varrho \tau_{j}^{*} \varrho\right)=\left(\varrho \tau_{j} \varrho\right)(\varrho \psi \varrho)$. By the special form (3.14) of $\varrho$ it is easy to see that $\varrho \psi \varrho$ also has the form (3.3) with the normalization (3.4). Also, $\varrho \tau_{j}^{*} \varrho$ is of the form (3.6). By the uniqueness part of Lemma 3.2 it follows that $\varrho \psi \varrho=\psi$ and consequently $\varrho \tau_{1}^{*} \varrho=\tau_{2}^{*}$. This gives the condition on $\Lambda_{1}$ and $\Lambda_{2}$.
Q.E.D.

ThEOREM 3.4. Let $\tau_{1}$ and $\tau_{2}$ be two involutions as in Lemma 3.2. Then there exists a transformation $\psi$ in $\mathbb{G}^{1}$ taking $\tau_{1}$ and $\tau_{2}$ into the form

$$
\begin{align*}
& \psi^{-1} \tau_{1} \psi:(\xi, \eta, \zeta) \mapsto\left(\Lambda \eta, \Lambda^{-1} \xi, \zeta\right) \\
& \psi^{-1} \tau_{2} \psi:(\xi, \eta, \zeta) \mapsto\left(\Lambda^{-1} \eta, \Lambda \xi, \zeta\right) \tag{3.15}
\end{align*}
$$

where $\Lambda=\lambda+\ldots, \operatorname{Re} \lambda>0$, is a formal power series in $\xi$ and the product $\xi \eta$. The most general transformation of the $\tau_{j}$ into this normal form is $\psi \circ \sigma$ where

$$
\begin{equation*}
\sigma:(\xi, \eta, \zeta) \mapsto(r(\xi \eta, \zeta) \xi, r(\xi \eta, \zeta) \eta ; f(\xi \eta, \zeta)) \tag{3.16}
\end{equation*}
$$

and $r(0,0) \neq 0$ and f is invertible. If in addition $\varrho \tau_{1}=\tau_{2} \varrho$, where $\varrho$ is given by (3.14), then $\psi \varrho=\varrho \psi$ and $r(\xi \eta, \zeta)=\bar{r}(\xi \eta, \zeta), f(\xi \eta, \xi)=\bar{f}(\xi \eta, \zeta) . \Lambda$ satisfies
(i) $\Lambda(\xi \eta, \zeta)=\bar{\Lambda}(\xi \eta, \zeta)$,
(ii) $\Lambda(\xi \eta, \zeta) \bar{\Lambda}(\xi \eta, \zeta)=1$,
according to the form of $\varrho$.

Proof. By the linear theory we may assume $\lambda_{1}=\lambda_{2}^{-1}=\lambda, \operatorname{Re} \lambda>0$. Consider the mapping

$$
(\xi, \eta, \zeta) \mapsto\left(v(\xi \eta, \zeta) \xi, v(\xi \eta, \zeta)^{-1} \eta, \zeta\right)
$$

which preserves $\xi \eta$. It commutes with $\varrho$ if
(i) $\nu \bar{\nu}=1$,
or
(ii) $\boldsymbol{v}=\bar{\nu}$.

Its effect is to preserve the form (3.6) while replacing $\Lambda_{j}$ by $\Lambda_{j} \nu^{-2}$. We can make $\Lambda_{1} \Lambda_{2}=1$ by choosing

$$
\begin{equation*}
v^{4}=\Lambda_{1} \Lambda_{2} \tag{3.17}
\end{equation*}
$$

If $\varrho$ is given by (3.14 i), then $\Lambda_{1} \bar{\Lambda}_{2}=1$, so there is a fourth root $v$ satisfying $v \bar{v}=1$. If $\varrho$ is given by ( 3.14 ii ), then $\Lambda_{1}=\bar{\Lambda}_{1}$ and there is a real fourth root. Let $\psi_{1} \in \mathcal{G}^{1}$ be a transformation of the $\tau_{j}$ into normal form $\tau_{j}^{*}$. We factor $\psi_{1}=\psi_{0} \circ \delta$ as in Lemma 3.1. It is easy to check that the transformation $\delta$ takes normalized involutions into normalized involutions. So $\psi_{0} \delta \tau_{j}^{*}=\tau_{j} \psi_{0} \delta$, or $\psi_{0}\left(\delta \tau_{j}^{*} \delta^{-1}\right)=\tau_{j} \psi_{0}$, implies $\psi_{0}=\psi$ and $\delta \tau_{j}^{*}=\tau_{j}^{*} \delta$ by the uniqueness statement in Lemma 3.2. This last relation gives

$$
\Lambda_{j}(\alpha \beta \xi \eta, \zeta+\gamma) \beta=\alpha \Lambda_{j}(\xi \eta, \zeta), \quad j=1,2 .
$$

Since $\Lambda_{1} \Lambda_{2}=1$, we get $\alpha^{2}=\beta^{2}$. The restriction $\operatorname{Re} \lambda>0$ forces $\alpha=\beta \equiv r$, and we set $f=\zeta+\gamma(\xi \eta, \zeta)$. If $\psi_{1} \notin \mathbb{F}^{1}$ its linear part has the form resulting from Lemma 2.3. This proves the theorem.
Q.E.D.

Let $Q_{\gamma}$ be one of the quadrics (2.1) with $\gamma \neq 0,1 / 2, \infty$. We shall say that $Q_{\gamma}$ is an exceptional hyperboloid if $\lambda$ given by (2.14) is a root of unity. Necessarily $\gamma>1 / 2$. If $\psi$ is a formal automorphism of $Q_{\gamma}$, it induces on $\mathcal{Q}_{\gamma}$ a mapping $\psi$ satisfying

$$
\psi \tau_{j}=\tau_{j} \psi, \quad \psi \varrho=\varrho \psi .
$$

By the theorem $\psi$ is of the form (3.16) with $r$ and $f$ real. Passing to $z, w$ coordinates via (2.11), (2.13) we get a mapping of the form (2.16). Hence, we have

Corollary 3.5. Suppose that $\gamma \neq 0,1 / 2, \infty$ and that $Q_{\gamma}$ is not an exceptional hyperboloid. Then the most general formal automorphism of $Q_{\gamma}$ is of the form (2.16).

The normal form for $\varphi$ lends itself to showing that $\varphi$ can be embedded in a flow $\varphi^{t}$ with $\varphi^{1}=\varphi, \varphi^{0}=\mathrm{id} ; \varphi^{t_{1}+t_{2}}=\varphi^{t_{1}} \circ \varphi^{t_{2}}$. We discuss this question for $n=2$, since the variables $\zeta$ are uninteresting for this problem. In the realm of formal power series a mapping

$$
\varphi:(\xi, \eta) \mapsto\left(\mu \xi+\ldots, \mu^{-1} \eta+\ldots\right)
$$

with $\mu$ not a root of unity can always be embedded in such a flow; moreover if $\varrho \varphi=\varphi^{-1} \varrho$ we have $\varrho \varphi^{t}=\varphi^{-i} \varrho$ and the embedding is essentially unique. The freedom is determined by the choice of $\log \mu$ alone.

The existence of such an interpolation follows at once from the normal form

$$
\varphi:(\xi, \eta) \mapsto\left(M \xi, M^{-1} \eta\right)
$$

by defining the formal power series

$$
N(\xi \eta)=\log \mu+\log \left(\mu^{-1} M\right)
$$

where the second term is a series without constant term. Thus the embedding is given by

$$
\varphi^{t}:(\xi, \eta) \mapsto\left(e^{t N} \xi, e^{-t N} \eta\right)
$$

If $\mu$ is real, and a posteriori positive, we have $M=\bar{M}$ and can define $N$ as a real series. If $|\mu|=1$ we have $M \bar{M}=1$ and we have in $N$ a series with purely imaginary coefficients, i.e. we have in the two cases
(i) $N=\bar{N}$,
(ii) $N+\bar{N}=0$.

This implies $\varrho \varphi^{t}=\varphi^{-i} \varrho$ in both cases.
This flow is generated by the $t$-independent vectorfield

$$
\dot{\xi}=N(\xi \eta) \xi, \quad \dot{\eta}=-N(\xi \eta) \eta
$$

which preserves the function $\xi \eta$.
We note that $\varphi^{t}$ is holomorphic if the transformation into the normal form converges. This fact will be of importance in section 5 in the description of the boundaries of analytic discs as orbits of these flows.

To establish that $\varphi^{t}$ is uniquely determined by $\varphi$ and the choice of $\log \mu$ we note that $\varphi^{t}$ has to commute with $\varphi$ and therefore is of the form

$$
\varphi^{t}:(\xi, \eta) \mapsto\left(\alpha_{r}(\xi \eta) \xi, \beta_{r}(\xi \eta) \eta\right)
$$

by our previous considerations. This corresponds to a differential equation of the form

$$
\begin{equation*}
\dot{\xi}=A(\xi \eta) \xi, \quad \dot{\eta}=B(\xi \eta) \eta \tag{}
\end{equation*}
$$

We claim that $A+B=0$. If this were not the case we would have

$$
A+B=c(\xi \eta)^{s}+\ldots, \quad c \neq 0
$$

for some $s \geqslant 1$ and therefore

$$
(\xi \eta)^{\cdot}=(A+B) \xi \eta=c(\xi \eta)^{s+1+\ldots}
$$

By integration of this formal differential equation we find

$$
(\xi \eta)(t)=(\xi \eta)+c t(\xi \eta)^{s+1}+\ldots
$$

For $t=1$ we see that $\varphi$ would not preserve $\xi \eta$, a contradiction. Hence $A+B=0$ and $\xi \eta$ is a constant for the differential equation $\left(^{*}\right)$ which can be integrated to

$$
\varphi^{t}:(\xi, \eta) \mapsto\left(e^{t A} \xi, e^{-t A} \eta\right)
$$

i.e. $e^{A}=M, A=\log M$, proving our claim.

There is another way to define $\varphi^{t}$ which goes back to G. D. Birkhoff [3]. If $\mu$ is not a root of unity one shows inductively that the iterates $\varphi^{j}, j=1,2, \ldots$ of $\varphi$ can be written in form

$$
\varphi^{j}:(\xi, \eta) \mapsto\left(\mu^{j} \xi+\sum_{d=2}^{\infty} f_{d}\left(\xi, \eta, \mu^{j}\right), \mu^{-j} \eta+\sum_{d=2}^{\infty} g_{d}\left(\xi, \eta, \mu^{j}\right)\right)
$$

where $f_{d}, g_{d}$, the homogeneous polynomials of degree $d$ in $\xi, \eta$, have coefficients which are polynomials in $\mu^{j}$ and $\mu^{-j}$ of degree $\leqslant d$. By replacing $\mu^{j}$ by $\mu^{t}=e^{t \log \mu}$ one obtains the formal series for $\varphi^{t}$. It was a fundamental observation of Birkhoff that-at least in the case of area preserving mappings-the series for $\varphi^{t}$ will in general diverge for noninteger $t$ even if $\varphi$ and hence $\varphi^{j}$ converges. In fact, in the case of area preserving mappings, the convergence of the transformation into normal form occurs precisely if this embedding can be achieved with convergent $\varphi^{t}$. The relation

$$
\varrho \circ \varphi^{t}=\varphi^{-i} \circ \varrho
$$

shows that $\varphi^{t}$ commutes with $\varrho$ precisely if $t+\bar{t}=0$, i.e. if $t$ is purely imaginary. This will imply that $\varphi^{t}$ for $t+\bar{t}=0$ gives rise to a flow on the real analytic manifold $M^{n}$ (see section 5).

## 4. Convergence

In general the transformation $\psi$ of Lemma 3.2 taking the pair of involutions $\tau_{1}, \tau_{2}$ into the normal form (3.6) does not converge even if $\tau_{1}$ and $\tau_{2}$ are given by convergent series. However, the following result gives a sufficient condition for convergence. It is proved by a majorant argument along the lines of the argument given in [12] and [11] for hyperbolic area preserving mappings.

THEOREM 4.1. Let the involutions $\tau_{1}, \tau_{2}$ be given by (3.1), where $p_{j}, q_{j}, r_{j}, j=1,2$, are convergent power series. If $\left|\lambda_{1}\right| \neq\left|\lambda_{2}\right|$ then the normalized transformation $\psi$ and the factors $\Lambda_{1}, \Lambda_{2}$ of Lemma 3.2 are given by convergent power series.

Proof. We make use of the following notations. Let $f(x), g(x), \ldots, h(y)$ be power series in some variables $x$ and $y$. If $f(x)=\Sigma a_{I} x^{I}$ (multi-index notation) then $f^{*}(x)=\Sigma\left|a_{I}\right| x^{I}$. Also, $f<g$ means that $g$ has non-negative coefficients $b_{I}$, and $\left|a_{I}\right| \leqslant b_{I}$, for all $I$. Note that if $f_{i}<g_{i}$ and $h_{1}<h_{2}$ then $h_{1}\left(f_{1}, f_{2}, \ldots\right)<h_{2}\left(g_{1}, g_{2}, \ldots\right)$, if $f_{i}, g_{i}$ have no constant terms.

Our argument will be based on the fact that $\psi$, given by (3.3, 3.4), transforms $\varphi$, given by (3.2), into the normal form $\tilde{\varphi}$, given in (3.6). The relation $\psi \circ \tilde{\varphi}=\varphi \circ \psi$ gives the functional equations

$$
\begin{align*}
& U\left(M \xi, M^{-1} \eta, \zeta\right)-\mu U(\xi, \eta, \zeta)=f(U, V, W) \\
& V\left(M \xi, M^{-1} \eta, \zeta\right)-\mu^{-1} V(\xi, \eta, \zeta)=g(U, V, W)  \tag{4.1}\\
& W_{a}\left(M \xi, M^{-1} \eta, \zeta\right)-W_{a}(\xi, \eta, \zeta)=h_{a}(U, V, W)
\end{align*}
$$

We decompose these equations by equating terms of the same type $s,-\infty<s<+\infty$, (see the definition before (3.4)),

$$
\begin{align*}
& \left(M^{s}-\mu\right) U_{s}=[f(U, V, W)]_{s} \\
& \left(M^{s}-\mu^{-1}\right) V_{s}=[g(U, V, W)]_{s}  \tag{4.2}\\
& \left(M^{s}-1\right)\left(W_{a}\right)_{s}=\left[h_{a}(U, V, W)\right]_{s} .
\end{align*}
$$

By interchanging $\tau_{1}$ and $\tau_{2}$ if necessary, we may assume that $\left|\lambda_{1}\right|>\left|\lambda_{2}\right|$, i.e. $|\mu|>1 . M^{-1}$ is a formal power series with constant term $\mu^{-1},\left|\mu^{-1}\right|<1$. Let $P^{0}=\left(M^{-1}-\mu^{-1}\right)^{*}$ and $P=|\mu|^{-1}+P^{0}$, so that $M^{-1}<P$.

We next prove the following relations.

$$
\begin{equation*}
\left(M^{s}-\mu^{k}\right)^{-1}<\frac{c}{1-c P^{0}}, \quad s \in \mathbf{Z}, k=0, \pm 1, s \neq k \tag{4.3}
\end{equation*}
$$

where the constant $c$ is independent of $s$. We first take $c \geqslant\left(1-\mu^{k}\right)^{-1}$ and assume $s \neq 0$. We also note that $|\mu|^{k / s} \leqslant \sqrt{|\mu|}$ whenever $s \neq k$ and $s \neq 0$. For $s \geqslant 1$, we have

$$
\begin{aligned}
\left(M^{s}-\mu^{k}\right)^{-1} & =M^{-s}\left(1-\mu^{k} M^{-s}\right)^{-1} \\
& =M^{-s} \sum_{j=0}^{\infty}\left(\mu^{k} M^{-s}\right)^{j} \\
& <P^{s} \sum\left(|\mu|^{k s} P\right)^{s j} \\
& <(\sqrt{|\mu|} P)^{s} \sum(\sqrt{|\mu|} P)^{s j} \\
& <(1-\sqrt{|\mu|} P)^{-1} .
\end{aligned}
$$

Thus (4.3) hold for $s \geqslant 1$ with $c \geqslant|\mu|(\sqrt{|\mu|}-1)^{-1}$. Now let $s=-t \leqslant-1$. Then

$$
\begin{aligned}
\left(M^{s}-\mu^{k}\right)^{-1} & =-\mu^{-k}\left(1-\mu^{-k} M^{s}\right)^{-1} \\
& =-\mu^{-k} \sum_{j=0}^{\infty}\left(\mu^{-k} M^{s}\right)^{j} \\
& <|\mu|^{-k} \sum\left(|\mu|^{-k t} P\right)^{i j} \\
& <|\mu| \sum_{a=0}^{\infty}(\sqrt{|\mu|} P)^{\alpha}=|\mu|(1-\sqrt{\mu} P)^{-1} .
\end{aligned}
$$

Hence, (4.3) holds for $s \leqslant-1$ if $c \geqslant|\mu|^{3 / 2}\left(|\mu|^{1 / 2}-1\right)^{-1}$.
From (4.2), (4.3), and (3.4) we get

$$
\begin{aligned}
U-\xi & =\sum_{s \neq 1} U_{s}=\sum_{s \neq+1}\left(M^{s}-\mu\right)^{-1}[f(U, V, W)]_{s} \\
& <\frac{c}{1-c P^{0}} \sum_{s}[f(U, V, W)]_{s}^{*} \\
& =\frac{c}{1-c P^{0}}[f(U, V, W)]^{*} .
\end{aligned}
$$

This gives the first of the three relations

$$
\begin{align*}
u & =U-\xi<\frac{c}{1-c P^{0}} f^{*}\left(U^{*}, V^{*}, W^{*}\right) \\
v & =V-\eta<\frac{c}{1-c P^{0}} g^{*}\left(U^{*}, V^{*}, W^{*}\right)  \tag{4.4}\\
w_{a} & =W_{a}-\zeta_{a}<\frac{c}{1-c P^{0}} h_{a}^{*}\left(U^{*}, V^{*}, W^{*}\right)
\end{align*}
$$

The second two are proved similarly. If we set $s=-1$ in the second equation of (4.2), we get

$$
\left(M^{-1}-\mu^{-1}\right) \eta=[g(U, V, W)]_{-1}
$$

from which follows

$$
\begin{equation*}
P^{0} \eta<[g(U, V, W)]^{*}<g^{*}\left(U^{*}, V^{*}, W^{*}\right) \tag{4.5}
\end{equation*}
$$

Since $f, g, h_{\alpha}$ converge and begin with quadratic terms we have a relation

$$
\begin{equation*}
f, g, h_{a}<G\left(x+y+\sum_{a} z_{a}\right), \quad G(t)=\frac{c_{1} t^{2}}{1-c_{1} t} \tag{4.6}
\end{equation*}
$$

for some $c_{1}>0$. Now we set $\eta=\zeta_{\alpha}=\xi$ and define the power series $W(\xi)$ by

$$
\xi W(\xi)=u^{*}+v^{*}+\sum_{\alpha=2}^{n-1} w_{\alpha}^{*}+P^{0} \xi
$$

Note that $W(0)=0$. From (4.4), (4.5), (4.6) we get

$$
\xi W<\left(\frac{n c}{1-c W}+1\right) G(Y)
$$

where

$$
\begin{aligned}
Y & =U^{*}+V^{*}+\sum_{a} W_{a}^{*} \\
& =n \xi+u^{*}+v^{*}+\sum w_{a}^{*} \\
& <\xi(n+W) .
\end{aligned}
$$

Hence,

$$
W<\frac{c_{2}}{1-c_{2} W} \frac{c_{1} \xi(n+W)^{2}}{1-c_{1} \xi(n+W)},
$$

for a suitable constant $c_{2}$. It follows that the series $W(\xi)$ is majorized by the solution $X=X(\xi), X(0)=0$, of the cubic equation

$$
X\left(1-c_{2} X\right)\left(1-c_{1} \xi(n+X)\right)=c_{1} c_{2} \xi(n+X)^{2}
$$

which is analytic near $\xi=0$. It follows that $u, v, w$ converge when $\xi=\eta=\zeta_{\alpha}$ have some non-zero value, and hence in a neighborhood of the origin. From the convergence of the map $\psi$ it follows that $\Lambda_{1}$ and $\Lambda_{2}$ converge.

It is clear that the $v$ given by (3.17) converges if $\tau_{1}$ and $\tau_{2}$ are given by convergent power series.

COROLLARY 4.2. Let the holomorphic involutions $\tau_{1}, \tau_{2}$ be as in the above theorem. Then the transformation $\psi$ and the factor $\Lambda$ in (3.15) are holomorphic.

## 5. Normal form for surfaces

In this section we use the results of section 3,4 on the normalization of the involutions $\tau_{1}, \tau_{2}, \varrho$ to transform the surface $M^{n} \subset C^{n}$ into a normal form near a suitable complex
tangent. Let $\left(\xi_{*}, \eta_{*}, \xi_{*}\right)$ be the normal coordinates on the complexified surface $\mathfrak{M}^{n} \subset \mathbf{C}^{2 n}$. Then we have

$$
\tau_{1}:\left\{\begin{array} { l } 
{ \xi _ { * } ^ { \prime } = \Lambda _ { * } \eta _ { * } } \\
{ \eta _ { * } ^ { \prime } = \Lambda _ { * } ^ { - 1 } \xi _ { * } , \quad \tau _ { 2 } : } \\
{ \zeta _ { * } ^ { \prime } = \zeta _ { * } }
\end{array} \quad \left\{\begin{array}{l}
\xi_{*}^{\prime}=\Lambda_{*}^{-1} \eta_{*} \\
\eta_{*}^{\prime}=\Lambda_{*} \zeta_{*} \\
\zeta_{*}^{\prime}=\zeta_{*}
\end{array}\right.\right.
$$

where $\Lambda_{*}=\Lambda_{*}\left(\xi_{*} \eta_{*}, \zeta_{*}\right)=\lambda+\ldots$ and in the two cases
(i) $\varrho\left(\xi_{*}, \eta_{*}, \zeta_{*}\right)=\left(\bar{\eta}_{*}, \bar{\xi}_{*}, \bar{\xi}_{*}\right), \quad \bar{\Lambda}_{*}=\Lambda_{*}$
or
(ii) $\varrho\left(\xi_{*}, \eta_{*}, \xi_{*}\right)=\left(\bar{\xi}_{*}, \bar{\eta}_{*}, \bar{\zeta}_{*}\right), \quad \bar{\Lambda}_{*} \Lambda_{*}=1$.

We still have the freedom to replace $\left(\xi_{*}, \eta_{*}, \zeta_{*}\right)$ by

$$
\begin{equation*}
\xi_{*}=r \xi, \quad \eta_{*}=r \eta, \quad \zeta_{*}=\xi \tag{5.1}
\end{equation*}
$$

$r=r(\xi \eta, \zeta)$, leading to $\Lambda(\xi \eta, \zeta)=\Lambda_{*}\left(\xi \eta r^{2}, \zeta\right)$.
Of course, also $\zeta_{*}$ can be reparametrized, but we will not make use of this fact, and determine $r$ in such a way that the surface is in a simple normal form.

ThEOREM 5.1. Assume that $M^{n}$ is a real analytic surface in $\mathbf{C}^{n}$ given by (1.16, 1.17) with $0<\gamma<1 / 2$, i.e. in the elliptic case. Then there exists a biholomorphic transformation near the origin taking $M^{n}$ into the implicit form

$$
\begin{align*}
& x_{n}=z_{1} \bar{z}_{1}+\Gamma\left(x_{n}, x_{\alpha}\right)\left(z_{1}^{2}+\bar{z}_{1}^{2}\right) \\
& y_{n}=0  \tag{5.2}\\
& y_{\alpha}=0 \quad(\alpha=2,3, \ldots n-1)
\end{align*}
$$

where $\Gamma=\bar{\Gamma}=\gamma+\ldots$.
Proof. As in the linear case (section 2) we introduce

$$
\begin{aligned}
& z_{1}=i \Lambda^{-1 / 2}(\Lambda \xi+\eta) \\
& w_{1}=-i \Lambda^{-1 / 2}(\xi+\Lambda \eta) \\
& z_{\alpha}=w_{\alpha}=x_{\alpha}=\zeta_{\alpha}
\end{aligned}
$$

These are equations on $\mathfrak{M}$. The third equation means that $z_{a}$ is the extension of $\zeta_{a}$ holomorphic in the original $z$ 's, and $w_{\alpha}$ is the extension of $\zeta_{\alpha}$ holomorphic in the
original $w$ 's. Hence, the ( $z_{\alpha}, w_{\alpha}$ ) are independent functions in the ambient space $\mathbf{C}^{2 n}$. The same computation as in section 2 gives

$$
z_{1} w_{1}+\Gamma_{1}(\xi \eta, \xi)\left(z_{1}^{2}+w_{1}^{2}\right)=\Delta(\xi \eta, \zeta)
$$

where

$$
\begin{aligned}
\Gamma_{1} & =\left(\Lambda+\Lambda^{-1}\right)^{-1} \\
\Delta & =\frac{\left(\Lambda-\Lambda^{-1}\right)^{2}}{\Lambda+\Lambda^{-1}} \xi \eta=\left(\Gamma_{1}^{-1}-4 \Gamma_{1}\right) \xi \eta
\end{aligned}
$$

We set

$$
z_{n}=w_{n}=\Delta(\xi \eta, \zeta)
$$

Note $z_{j}$ are $\tau_{2}$-invariant, while $w_{j}$ are $\tau_{1}$-invariant and therefore they are holomorphically or antiholomorphically related to the original coordinates $z_{j}^{0}$. Because of the choice of the linear terms we have in

$$
\begin{aligned}
z_{j} & =f_{j}\left(z^{0}\right)=z_{j}^{0}+\ldots \\
w_{j} & =g_{j}\left(w^{0}\right)=w_{j}^{0}+\ldots
\end{aligned}
$$

a biholomorphic change of coordinates in $\mathbf{C}^{2 n}$. Because of the reality conditions we have $w_{1}=c z_{1} 0 \varrho$ and $g_{j}\left(z^{0}\right)=f_{j}\left(z^{0}\right)$. Hence $z^{0} \mapsto z$ defines a holomorphic coordinate change in $\mathrm{C}^{n}$.

Now if we eliminate $\xi \eta, \zeta$ from

$$
x_{n}=\Delta(\xi \eta, \zeta) ; \quad x_{\alpha}=\zeta_{\alpha}
$$

and set

$$
\Gamma\left(x_{n}, x_{a}\right)=\Gamma_{1}(\xi \eta, \zeta)
$$

we obtain (5.2).
Q.E.D.

In the hyperbolic case, if $\mu$ is not a root of unity we find the same normal form if we admit coordinate transformations given only by formal series, since, in general, we have to expect divergence.

Assuming convergence we read off several important facts about $M$. First $M$ admits the holomorphic involution $\left(z_{1}, z_{\alpha}, z_{n}\right) \mapsto\left(-z_{1}, z_{\alpha}, z_{n}\right)$. Also, $M$ lies in the linear space $\operatorname{Im} z_{\alpha}=\operatorname{Im} z_{n}=0$. An $(n-1)$-real parameter family of complex lines cut $M$ in a family of disjoint real analytic curves. In the elliptic case, where convergence is
guaranteed, these are closed real curves bounding linear analytic discs on $z_{\alpha}=c_{\alpha}=\bar{c}_{\alpha}$, $z_{n}=c_{n}=\bar{c}_{n}>0$, bounded by the ellipses

$$
c_{n}=z_{1} \bar{z}_{1}+\Gamma\left(c_{n}, c_{a}\right)\left(z_{1}^{2}+\bar{z}_{1}^{2}\right)
$$

These discs sweep out a $(n+1)$-dimensional real analytic manifold $\tilde{M}$ with boundary $M$. This $\tilde{M}$ is the local holomorphic hull of $M$. There are no other analytic discs $\mathfrak{C}$ in $\mathbf{C}^{n}$ with boundaries on $M$ near 0 . Indeed, the functions $\operatorname{Im} z_{\alpha}, \operatorname{Im} z_{n}$ vanish on the boundary of $\mathfrak{C}$, hence identically on $\mathfrak{C}$. Thus $z_{\alpha}, z_{n}$ are real constants on $\mathfrak{C}$, and so $\mathfrak{C}$ lies on the discs given by $z_{\alpha}=c_{\alpha}, z_{n}=c_{n}$.

In the hyperbolic case, provided we have a convergent transformation $\psi$ into normal form, this argument shows that there exists no analytic disc with boundary on $M^{n}$ near 0 .

Next we make use of the transformation (5.1) to further simplify the factor $\Gamma$. We distinguish two cases: In the first $\Gamma_{*}$ is independent of $\xi_{*} \eta_{*}$, hence $\Gamma$ is independent of $x_{n}$; in this case no further simplification is achieved. In the second case we write

$$
\Gamma_{*}\left(\xi_{*} \eta_{*}, \zeta_{*}\right)=\gamma\left(\zeta_{*}\right)+\sum_{k>s} \gamma_{k}\left(\zeta_{*}\right)\left(\xi_{*} \eta_{*}\right)^{k}
$$

where $\gamma_{s}\left(\zeta_{*}\right) \neq 0 ; \gamma(0)=\gamma$. The integer $s$ is a biholomorphic invariant. Since there are points $\zeta_{*}$ with $\gamma\left(\zeta_{*}\right) \neq 0$ near 0 we may assume that $\gamma_{s}(0) \neq 0$ and will choose $r=n\left(\xi_{*} \eta_{*}, \zeta_{*}\right)$ in (5.1) so that $\Gamma$ has the form

$$
\begin{equation*}
\Gamma=\gamma\left(x_{\alpha}\right)+\delta x_{n}^{s}, \quad \delta= \pm 1 . \tag{5.3}
\end{equation*}
$$

Indeed, since $\Gamma\left(x_{a}, x_{n}\right)$ is obtained by elimination of $\xi \eta$ from

$$
\begin{gathered}
x_{n}=\Delta(\xi \eta, \xi)=\left(\Gamma_{1}^{-1}-4 \Gamma_{1}\right) \xi \eta \\
\Gamma\left(x_{a}, x_{n}\right)=\Gamma_{1}\left(\xi \eta, x_{a}\right)=\Gamma_{*}\left(\xi \eta r^{2}, x_{a}\right)
\end{gathered}
$$

we have to solve the equation

$$
\Gamma_{*}\left(\xi \eta r^{2}, x_{\alpha}\right)=\gamma\left(x_{\alpha}\right)+\delta \Delta^{s}
$$

or with $\zeta_{\alpha}=x_{\alpha}$

$$
\left[\Gamma_{*}^{-1}\left(\xi \eta r^{2}, \zeta\right)-4 \Gamma_{*}\left(\xi \eta r^{2}, \zeta\right)\right]^{s}=\delta^{-1}\left(\Gamma_{*}\left(\xi \eta r^{2}, \zeta\right)-\gamma(\xi)\right)=\delta^{-1} r^{2 s} \sum_{j>0} \gamma_{j+s} r^{2 j}(\xi \eta)^{j}
$$

We choose $\delta= \pm 1$ so that $\delta \gamma^{-s}\left(1-4 \gamma^{2}\right)^{s} \gamma_{s}(0)^{-1}>0$; recall that the coefficients of $\Gamma_{*}$ are real. Taking the $(2 s)^{\text {th }}$ root of both sides we can solve for a real $r$ by the implicit function theorem. The function $r$ is determined up to sign. Thus we have achieved the form (5.3) for $\Gamma$.

We may still apply a real invertible transformation $\zeta \mapsto f(\zeta) ; f(\zeta)=\vec{f}(\zeta), f(0)=0$ to simplify $\gamma(\zeta)$. For example, if $\zeta=0$ is a regular point we could achieve $\gamma(\xi)=\gamma+\xi_{2}$, or $\gamma\left(x_{\alpha}\right)=\gamma+x_{2}$. Thus, generically $M^{n}$ has the form

$$
\begin{align*}
& x_{n}=z_{1} \bar{z}_{1}+\left(\gamma+x_{2} \pm x_{n}^{s}\right)\left(z_{1}^{2}+\bar{z}_{1}^{2}\right), \quad y_{a}=y_{n}=0, \quad \text { for } n \geqslant 3 \\
& x_{2}=z_{1} \bar{z}_{1}+\left(\gamma \pm x_{2}^{s}\right)\left(z_{1}^{2}+\bar{z}_{1}^{2}\right), \quad y_{2}=0, \quad \text { for } n=2 . \tag{5.4}
\end{align*}
$$

For $n=2$ the automorphism group of $M^{2}$ consists only of $\left(z_{1}, z_{2}\right) \mapsto\left( \pm z_{1}, z_{2}\right)$ provided that $\lambda$ is not a root of unity and $\gamma_{s} \neq 0$, i.e. in case that the normal form does not represent a quadric. This follows readily from our formal considerations.

We recall that the involutions $\tau_{j}$ map $\mathfrak{M}$ into itself, hence also $\varphi$ and $\varphi^{t}$, defined in section 3 map $\mathfrak{M}$ into itself. However, in order that $\varphi^{t}$ maps $M$, the fixed point set of $\varrho$, into itself, we need that $\varrho$ and $\varphi^{t}$ commute. As was shown at the end of section 3 this is the case for purely imaginary $t=i \sigma$. In the elliptic case the orbits of $\varphi^{i \sigma}, \sigma$ real, are the closed curves which bound the analytic discs.

It suffices to prove this in the normal form. Since $\varphi^{i \sigma}$ preserves $\xi \eta$ as well as $\xi_{\alpha}$ it follows that the orbits lie on

$$
\zeta_{\alpha}=c_{\alpha} ; \quad \xi \eta=c_{n}
$$

where $c_{\alpha}, c_{n}$ are real constants, $c_{n}>0$. Hence we have

$$
z_{\alpha}=c_{\alpha}=w_{\alpha} ; \quad z_{n}=\text { constant }
$$

which proves the claim.

## 6. Further remarks

(a) Exceptional hyperbolic surfaces. We consider a surface $M$ in $\mathbf{C}^{2}$ given by (1.3) with $\gamma>1 / 2, H=h+i k$. We assume that the associated mapping $\varphi$ on $\mathfrak{M}$ is such that $\varphi^{\prime}(0)$ is nilpotent,

$$
\lambda^{2 m}=1, \lambda^{2 j} \neq 1, \quad j<m
$$

We shall show that $M$ can be holomorphically flattenend to order $m$ and in general to no higher order. For this it suffices to consider transformations of the form

$$
\begin{equation*}
\tilde{z}_{1}=z_{1}, \quad \tilde{z}_{2}=z_{2}+B\left(z_{1}, z_{2}\right) \tag{6.1}
\end{equation*}
$$

where $B\left(z_{1}, z_{2}\right)$ is polynomial without constant or linear terms. Restriction to $M$ yields

$$
\tilde{z}_{2}=q+H+B\left(z_{1}, q+H\right) \equiv q+\tilde{H}
$$

so that $\operatorname{Im} B\left(z_{1}, q+H\right)=\tilde{k}-k$. Suppose that $k$ and $\tilde{k}$ begin with terms of degree $n<m$. We shall choose a holomorphic polynomial $B\left(z_{1}, z_{2}\right)$ of weight $n$ (weight of $z_{j}=j, j=1,2$ ) so as to annihilate the terms of degree $n$ in $\tilde{k}$. This amounts to solving an equation of the form

$$
\begin{equation*}
\operatorname{Im} B\left(z_{1}, q\right)=k \tag{6.2}
\end{equation*}
$$

where $k$ is a real homogeneous polynomial of degree $n$. This is a problem on the quadric $Q_{\gamma}$.

Complexifying gives

$$
\frac{1}{2 i} B\left(z_{1}, q\left(z_{1}, w_{1}\right)\right)-\frac{1}{2 i} \tilde{B}\left(w_{1}, q\left(z_{1}, w_{1}\right)\right)=k\left(z_{1}, w_{1}\right)
$$

which implies that $k$ can be decomposed into the sum of a homogeneous polynomial of degree $n$ invariant under $\tau_{2}$ and one invariant under $\tau_{1}$. We pass to the ( $\xi, \eta$ )-coordinate system by the linear change (2.11). The most general such polynomials invariant under $\tau_{1}$ and $\tau_{2}$ are

$$
f_{1}=\sum_{j \leqslant n / 2} a_{j}\left(\xi^{n-j} \eta^{j}+\lambda^{n-2 j} \xi^{j} \eta^{n-j}\right)
$$

and

$$
f_{2}=\sum_{j \leqslant n / 2} b_{j}\left(\xi^{n-j} \eta^{j}+\lambda^{2 j-n} \xi^{j} \eta^{n-j}\right)
$$

respectively. The real polynomial $k$ has the form

$$
k=\sum_{j=0}^{n} c_{j} \xi^{n-j} \eta^{j}, \quad c_{j}=\bar{c}_{j}
$$

so that $f_{1}+f_{2}=k$ reduces to

$$
\begin{align*}
a_{j}+b_{j} & =c_{j}, \quad 0 \leqslant j \leqslant n / 2,  \tag{6.3}\\
\lambda^{n-2 j} a_{j}+\lambda^{2 j-n} b_{j} & =c_{n-j}, \quad 0 \leqslant j \leqslant n / 2 .
\end{align*}
$$

For $2 j<n$ these equations have a unique solution $a_{j}, b_{j}$ since the determinant does not vanish because of $\lambda^{2(n-2 j} \neq 1$. Conjugating (6.3) we find $b_{j}=\bar{a}_{j}$ since $c_{j}=\bar{c}_{j}, \lambda \bar{\lambda}=1$. For $2 j=n$ the two equations agree with the single equation

$$
a_{j}+b_{j}=c_{j} .
$$

We choose the solution

$$
a_{j}=b_{j}=\frac{1}{2} c_{j}
$$

so that again $b_{j}=\bar{a}_{j}$ holds, since $c_{j}$ is real. Thus $f_{1}=\bar{f}_{2}$ is the trace of a function $(1 / 2 i) B$ holomorphic in $z_{1}, z_{2}$ which gives the solution of (6.2).

The first instance in which (6.3) may not be solvable is $n=m, j=0$. If $\lambda^{m}=-1$, we have the compatibility condition $c_{m}=-c_{0}$. For example, $k=\overline{k \circ \varrho}=\xi^{m}+\eta^{m}$ cannot be written as such a sum $f_{1}+f_{2}$. We set $\lambda_{m}=e^{i \pi / m}$, and by (2.14) $\gamma_{m}=(1 / 2) \sec (\pi / m)$. In particular $\gamma_{3}=1, \gamma_{4}=1 / \sqrt{2}$. The corresponding surface is (via (2.11))

$$
\begin{equation*}
z_{2}=\gamma_{m} z_{1}^{2}+z_{1} \bar{z}_{1}+\gamma_{m} \bar{z}_{1}^{2}+i^{m+1}\left(\lambda_{m}\right)^{m / 2}\left(\lambda_{m}^{2}-1\right)^{-m}\left\{\left(-\lambda_{m} z_{1}-\bar{z}_{1}\right)^{m}+\left(z_{1}+\lambda_{m} \bar{z}_{1}\right)^{m}\right\} \tag{6.4}
\end{equation*}
$$

The imaginary part of the right hand side cannot be made to vanish to higher order.
If $\lambda^{m}=+1$, we may take $k=\xi^{m}-\eta^{m}$.
For $\gamma=\gamma_{3}=1$, a simpler example of a surface which cannot be flattened to third order is

$$
\begin{equation*}
z_{2}=z_{1}^{2}+z_{1} \bar{z}_{1}+\bar{z}_{1}^{2}+z_{1} \bar{z}_{1}\left(z_{1}-\bar{z}_{1}\right) \tag{6.5}
\end{equation*}
$$

This can be seen by examining (6.2) directly in $\left(z_{1}, \tilde{z}_{1}\right)$-coordinates.
(b) Divergence in the normal form. When $\gamma>1 / 2$ is not exceptional, the results of section 5 show that the surface $M$ can be formally transformed into a real hyperplane. However, as mentioned before the transformation will in general diverge. Rather than
prove a general theorem to this effect, we shall give an example of a surface in $\mathbf{C}^{2}$ which cannot be holomorphically flattened. This surface will be of the form

$$
M: \begin{align*}
& z_{2}=\left(k\left(z_{1}\right)+\gamma \bar{z}_{1}\right) \bar{z}_{1}  \tag{6.6}\\
& k\left(z_{1}\right)=z_{1}+k_{0}\left(z_{1}\right),
\end{align*}
$$

where $k_{0}$ is a holomorphic polynomial in $z_{1}$ beginning with a term of order $\geqslant 2$.
If $M$ could be transformed into the hyperplane $\operatorname{Im} \tilde{z}_{2}=0$, then this could be accomplished by means of a transformation of the form (6.1) with $B=\gamma z_{1}^{2}+O\left(|z|^{3}\right)$. Let $G\left(z_{1}, \bar{z}_{1}\right)=q\left(z_{1}, \bar{z}_{1}\right)+\ldots$ be the restriction of $z_{2}+B$ to $M$. Then $G\left(z_{1}, \bar{z}_{1}\right)$ is also the restriction to $M$ of $\bar{z}_{2}+\bar{B}$. Consequently, the complex function $G=G\left(z_{1}, w_{1}\right)$ on $\mathfrak{M}$ is invariant under both $\tau_{1}$ and $\tau_{2}$, and so $G \circ \varphi=G$. Furthermore,

$$
d G=q_{z_{1}} d z_{1}+q_{w_{1}} d w_{1}+\ldots
$$

is non-zero in a deleted neighborhood of $z_{1}=w_{1}=0$. We shall show that if $1 / 2<\gamma<\infty$, $\gamma \neq 1 / \sqrt{2}$, then $k$ can be chosen so that $\varrho$ admits no such (non-trivial) invariant function $G$ in any neighborhood of the origin.

One readily sees that

$$
\begin{gather*}
\tau_{1}:\left(z_{1}, w_{1}\right) \mapsto\left(-z_{1}-\gamma^{-1} k\left(w_{1}\right), w_{1}\right), \\
\tau_{2}:\left(z_{1}, w_{1}\right) \mapsto\left(z_{1},-w_{1}-\gamma^{-1} k\left(z_{1}\right)\right), \\
\varphi: \begin{array}{l}
z_{1}^{\prime}=-z_{1}-\gamma^{-1} k\left(w_{1}^{\prime}\right) \\
w_{1}^{\prime}=-w_{1}-\gamma^{-1} k\left(z_{1}\right)
\end{array}, \quad \varphi^{-1}: \begin{array}{l}
z_{1}=-z_{1}^{\prime}-\gamma^{-1} k\left(w_{1}^{\prime}\right) \\
w_{1}=-w_{1}^{\prime}-\gamma^{-1} k\left(z_{1}\right)
\end{array} . \tag{6.7}
\end{gather*}
$$

In particular both $\varphi$ and $\varphi^{-1}$ are polynomial mappings (so called Cremona transformations) of the form

$$
\varphi:\left\{\begin{array}{l}
z_{1}^{\prime}=c z_{1}^{2 \delta}+\ldots \\
w_{1}^{\prime}=0+\ldots
\end{array}, \quad \varphi^{-1}:\left\{\begin{array}{l}
z_{1}=0+\ldots \\
\omega_{1}=\bar{c} w_{1}^{2 d}+\ldots
\end{array}\right.\right.
$$

where $\delta=\operatorname{deg} k$ and the dots indicate terms of lower degree in $\left(z_{1}, w_{1}\right)$ or $\left(z_{1}^{\prime}, w_{1}^{\prime}\right)$. The $n$-fold iterates of $\varphi$ and $\varphi^{-1}$ are of the form

$$
\varphi^{n}:\left\{\begin{array}{l}
z_{1}^{\prime}=c_{n} z_{1}^{(2 \delta)^{n}}+\ldots, \quad \varphi^{-n}:\left\{\begin{array}{l}
z_{1}=0+\ldots \\
w_{1}^{\prime}=0+\ldots
\end{array}, \bar{c}_{n} w_{1}^{(2 \delta)^{n}}+\ldots\right.
\end{array}\right.
$$

where $c_{n} \neq 0$. A fixed point $p$ of $\varphi^{2 n}$ satisfies $\varphi^{n}(p)=\varphi^{-n}(p)$, so is a solution to the pair of polynomial equation

$$
c_{n} z_{1}^{(2 \delta)^{n}}+\ldots=\bar{c}_{n} w_{1}^{(2 \delta)^{n}}+\ldots=0 .
$$

The leading terms show that these two polynomials have no common factor. By Bezout's theorem $\varphi^{2 n}$ can have at most ( $\left.2 \delta\right)^{2 n}$ fixed points.

From (6.7) we see that $d z_{1}^{\prime} \wedge d w_{1}^{\prime}=-d z_{1} \wedge d w_{1}^{\prime}=d z_{1} \wedge d w_{1}$, so that the Jacobian determinant of $\varphi$ is identically one. Consider the holomorphic vector field

$$
X_{G}=G_{w_{1}} \frac{\partial}{\partial z_{1}}-G_{z_{1}} \frac{\partial}{\partial w_{1}} .
$$

If $G \circ \varphi=G$, then $X_{G}$ is invariant under $\varphi: d \varphi\left(X_{G}\right)=X_{G}$. This follows from the chain rule and the fact that $\varphi$ is area-preserving. Therefore, if $p$ is a fixed point of $\varphi^{2 n}$ then so is every point on the orbit through $p$ of the flow $\exp t X_{G}$. If $d G(p) \neq 0$, then this orbit is locally a smooth holomorphic curve. $\varphi^{2 n}$ would then have a continuum of fixed points, which is impossible. It follows that if we can choose $k\left(z_{1}\right)$ so that every deleted neighborhood of the origin contains a fixed point of $\varphi^{n}$, for some $n$, then $M$ cannot be holomorphically flattened.

To achieve this last property we shall appeal to the Birkhoff fixed point theorem. This theorem applies to area-preserving transformations of the real plane. We must therefore choose a $k$ with real coefficients so that $\varphi$ leaves invariant the plane of real $z_{1}, w_{1}$. We shall, in fact, take $k\left(z_{1}\right)=z_{1}+\gamma z_{1}^{3}+\ldots$ a polynomial with real coefficients so that

$$
\varphi: \begin{aligned}
& z_{1}^{\prime}=\left(\gamma^{-2}-1\right) z_{1}+\gamma^{-1} w+f+O(4) \\
& w_{1}^{\prime}=-\gamma^{-1} z_{1}-w_{1}+g,
\end{aligned}
$$

where

$$
f=\gamma^{-1} z_{1}^{3}+\left(w_{1}+\gamma^{-1} z_{1}\right)^{3}, \quad g=-z_{1}^{3} .
$$

If we subject this to the coordinate change (2.11), we get

$$
\varphi: \begin{aligned}
& \xi^{\prime}=\lambda^{2} \xi-i \sqrt{\lambda}\left(\lambda^{2}-1\right)^{-1}(\lambda f+g)+O(4) \\
& \eta^{\prime}=\lambda^{-2} \eta+i \sqrt{\lambda}\left(\lambda^{2}-1\right)^{-1}(f+\lambda g)+O(4),
\end{aligned}
$$

where

$$
\begin{align*}
& \lambda f+g=6 \lambda^{4}(i / \sqrt{\lambda})^{3} \xi^{2} \eta+\ldots \\
& f+\lambda g=6(i / \sqrt{\lambda})^{3} \xi \eta^{2}+\ldots \tag{6.8}
\end{align*}
$$

Here the dots indicate the other cubic terms. If $\lambda^{2} \neq \pm 1, \pm i$, i.e. $\gamma \neq 1 / 2,1 / \sqrt{2}, \infty$, then these other terms can be removed by a further coordinate change which does not alter
the two coefficients shown in (6.8). (See [11] § 23 (18), p. 158.) Hence, $\varphi$ can be transformed into

$$
\varphi: \begin{aligned}
& \xi^{\prime}=\lambda^{2}(1+i a \xi \eta) \xi+O(4) \\
& \eta^{\prime}=\lambda^{-2}(1-i a \xi \eta) \eta+O(4)
\end{aligned}, \quad a=\frac{6 \gamma}{\sqrt{4 \gamma^{2}-1}}>0
$$

Actually, the case $\lambda^{2}= \pm i$ i.e. $\gamma=1 / \sqrt{2}$ does not have to be excluded since the relevant terms $\eta^{3}, \xi^{3}$ in $\lambda f+g, f+\lambda g$, respectively, have zero coefficients.

By Birkhoff's theorem [11], p. 174, for each sufficiently small neighborhood $U$ of the origin in the plane of real $z_{1}, w_{1}$ there exists an integer $n$ such that $\varphi^{2 n}$ fixes a point different from $O$ in $U$. In particular, we have proved the following proposition.

Proposition 6.1. If $1 / 2<\gamma<\infty$, then the hyperbolic surface

$$
z_{2}=z_{1} \bar{z}_{1}+\gamma \bar{z}_{1}^{2}+\gamma z_{1}^{3} \bar{z}_{1}
$$

cannot be transformed into a real hyperplane by means of a (convergent) biholomorphic transformation.

This shows that the formal transformations, which for $\lambda$ not a root of unity, $|\lambda|=1$, exist, must be divergent. This example shows also that divergence can not be avoided by inequalities of the type $\left|\lambda^{j}-1\right| \geqslant c|j|^{-\nu}$ for all $j \geqslant 1$. Incidentally, the periodic orbits of $\varphi$, as well as its invariant curves, do not have any geometrically significance since they do not lie on $M$ but only on its complexification $\mathfrak{M}$.
(c) The case $\gamma=0$. Here one can apply the formal theory of [5], section 2. We state without proof some of the results of this theory for surfaces in $\mathbf{C}^{2} . M$ has the form

$$
M: \begin{aligned}
& z_{2}=z_{1} \bar{z}_{1}+\operatorname{Re} h\left(z_{1}\right)+z_{1} \bar{z}_{1} H\left(z_{1}, \bar{z}_{1}\right) \\
& h\left(z_{1}\right)=\sum_{j \geqslant k} c_{j} z_{1}^{j}, \quad k=3,4,5, \ldots, \infty .
\end{aligned}
$$

The integer $k$, which is the degree of the lowest pure $\bar{z}_{1}$-term, is a biholomorphic invariant. $M$ is formally equivalent to the quadric $Q_{0}: z_{2}=z_{1} \bar{z}_{1}$ if and only if $k=\infty$. Otherwise, $M$ may be formally transformed into the form

$$
z_{2}=z_{1} \bar{z}_{1}+z_{1}^{k}+\bar{z}_{1}^{k}+\operatorname{Re} \sum_{j>k} a_{j} z^{j}
$$

The complex numbers $a_{j}$ are not absolute invariants, since this normal form is still
subject to the action of the (formal) automorphism group of $Q_{0}$. This group can be shown to be made up precisely of those transformations of the form

$$
\begin{aligned}
& z_{1}^{\prime}=G\left(x_{2}\right) \frac{z_{1}-x_{2} b\left(x_{2}\right)}{1-\bar{b}\left(x_{2}\right) z_{1}} \\
& x_{2}^{\prime}=G\left(x_{2}\right) \hat{G}\left(x_{2}\right) x_{2}
\end{aligned}
$$

where $G$ and $b$ are arbitrary complex formal power series with $G(0) \neq 0$. There are probably infinitely many real valued invariants in the case $\gamma=0$.

One still has the projections $\pi_{1}$ and $\pi_{2}$ on the complexified surface $\mathfrak{M}$. The case in which $M$ is formally equivalent to $Q_{0}$ is characterized by $\pi_{1}$ (or $\pi_{2}$ ) being locally one-to-one except for collapsing an analytic curve to a point. Otherwise ( $k<\infty$ ) each $\pi_{i}$ is a $k$-fold branched covering.

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Received June 1, 1982


[^0]:    ${ }^{(1)}$ Alfred P. Sloan Fellow. Partially supported by NSF, Grant No. MCS 8100793.

