# Lempert theorem for strongly linearly convex domains 

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#### Abstract

In 1984 L. Lempert showed that the Lempert function and the Carathéodory distance coincide on non-planar bounded strongly linearly convex domains with realanalytic boundaries. Following his paper, we present a slightly modified and more detailed version of the proof. Moreover, the Lempert Theorem is proved for non-planar bounded strongly linearly convex domains.


The aim of this paper is to present a detailed version of the proof of the Lempert Theorem in the case of non-planar bounded strongly linearly convex domains. Lempert's original proof is presented only in the proceedings of a conference (see [6]) with a very limited access and at some places it was quite sketchy. We were encouraged by some colleagues to prepare an extended version of the proof, in which all doubts could be removed and some details of the proofs could be simplified. We hope to have done it below. Certainly, the idea of the proof belongs entirely to Lempert. The main differences we would like to draw attention to are:

- the results are obtained in the $\mathcal{C}^{2}$-smooth case;
- the notions of stationary mappings and $E$-mappings are separated;
- a geometry of domains is investigated only in neighborhoods of boundaries of stationary mappings (viewed as boundaries of analytic discs)this allows us to obtain localization properties for stationary mappings.

Additional motivation for presenting the proof is the fact, showed recently in [7, that the so-called symmetrized bidisc may be exhausted by strongly linearly convex domains. On the other hand, it cannot be exhausted by domains biholomorphic to convex ones ([1). Therefore, the equality of the Lempert function and the Carathéodory distance for strongly linearly convex domains does not follow directly from [5].

[^0]1. Introduction and results. Let us recall the objects we will deal with. Throughout the paper, $\mathbb{D}$ denotes the unit open disc on the complex plane, $\mathbb{T}$ is the unit circle and $p$ is the Poincare distance on $\mathbb{D}$.

Let $D \subset \mathbb{C}^{n}$ be a domain and let $z, w \in D, v \in \mathbb{C}^{n}$. The Lempert function is defined as

$$
\begin{equation*}
\widetilde{k}_{D}(z, w):=\inf \{p(0, \xi): \xi \in[0,1) \text { and } \exists f \in \mathcal{O}(\mathbb{D}, D): f(0)=z, f(\xi)=w\} \tag{1}
\end{equation*}
$$

The Kobayashi-Royden (pseudo)metric is

$$
\begin{equation*}
\kappa_{D}(z ; v):=\inf \left\{\lambda^{-1}: \lambda>0 \text { and } \exists f \in \mathcal{O}(\mathbb{D}, D): f(0)=z, f^{\prime}(0)=\lambda v\right\} . \tag{2}
\end{equation*}
$$

Note that

$$
\begin{align*}
& \widetilde{k}_{D}(z, w)=\inf \{p(\zeta, \xi): \zeta, \xi \in \mathbb{D} \text { and }  \tag{3}\\
&\exists f \in \mathcal{O}(\mathbb{D}, D): f(\zeta)=z, f(\xi)=w\}, \\
& \kappa_{D}(z ; v)=\inf \left\{|\lambda|^{-1} /\left(1-|\zeta|^{2}\right): \lambda \in \mathbb{C}_{*}, \zeta \in \mathbb{D}\right. \text { and }  \tag{4}\\
&\left.\exists f \in \mathcal{O}(\mathbb{D}, D): f(\zeta)=z, f^{\prime}(\zeta)=\lambda v\right\} .
\end{align*}
$$

If $z \neq w$ (resp. $v \neq 0$ ), a mapping $f$ for which the infimum in (3) (resp. in (44) is attained is called a $\widetilde{k}_{D}$-extremal (or a Lempert extremal) for $z, w\left(\right.$ resp. a $\kappa_{D}$-extremal for $\left.z, v\right)$. A mapping which is a $\widetilde{k}_{D}$-extremal or a $\kappa_{D}$-extremal will be called just an extremal or an extremal mapping.

We shall say that $f: \mathbb{D} \rightarrow D$ is a unique $\widetilde{k}_{D}$-extremal for $z, w$ (resp. a unique $\kappa_{D}$-extremal for $z, v$ ) if any other $\widetilde{k}_{D}$-extremal $g: \mathbb{D} \rightarrow D$ for $z, w$ (resp. $\kappa_{D}$-extremal for $z, v$ ) satisfies $g=f \circ a$ for some Möbius function $a$.

The next objects we are dealing with are the Carathéodory (pseudo) distance

$$
c_{D}(z, w):=\sup \{p(F(z), F(w)): F \in \mathcal{O}(D, \mathbb{D})\}
$$

and the Carathéodory-Reiffen (pseudo)metric

$$
\gamma_{D}(z ; v):=\sup \left\{\left|F^{\prime}(z) v\right|: F \in \mathcal{O}(D, \mathbb{D}), F(z)=0\right\}
$$

A holomorphic mapping $f: \mathbb{D} \rightarrow D$ is said to be a complex geodesic if $c_{D}(f(\zeta), f(\xi))=p(\zeta, \xi)$ for any $\zeta, \xi \in \mathbb{D}$.

Here is some more notation. Let $z_{1}, \ldots, z_{n}$ be the standard complex coordinates in $\mathbb{C}^{n}$ and $x_{1}, \ldots, x_{2 n}$ the corresponding real coordinates in $\mathbb{C}^{n}=\mathbb{R}^{n}+i \mathbb{R}^{n} \simeq \mathbb{R}^{2 n}$. We use $T_{D}^{\mathbb{R}}(a), T_{D}^{\mathbb{C}}(a)$ to denote the real and the complex tangent space to a $\mathcal{C}^{1}$-smooth domain $D$ at a point $a \in \partial D$, i.e. the sets

$$
\begin{aligned}
T_{D}^{\mathbb{R}}(a) & :=\left\{X \in \mathbb{C}^{n}: \operatorname{Re} \sum_{j=1}^{n} \frac{\partial r}{\partial z_{j}}(a) X_{j}=0\right\} \\
T_{D}^{\mathbb{C}}(a) & :=\left\{X \in \mathbb{C}^{n}: \sum_{j=1}^{n} \frac{\partial r}{\partial z_{j}}(a) X_{j}=0\right\}
\end{aligned}
$$

where $r$ is a defining function of $D$. Let $\nu_{D}(a)$ be the outward unit normal vector to $\partial D$ at $a$.

Let $\mathcal{C}^{k}(\overline{\mathbb{D}})$, where $k \in(0, \infty]$, denote the class of continuous functions on $\overline{\mathbb{D}}$ which are of class $\mathcal{C}^{k}$ on $\mathbb{D}$ and

- if $k \in \mathbb{N} \cup\{\infty\}$ then derivatives in $\mathbb{D}$ up to order $k$ extend continuously on $\overline{\mathbb{D}}$;
- if $k-[k]=: c>0$ then derivatives in $\mathbb{D}$ up to order $[k]$ extend to $c$-Hölder continuous functions on $\overline{\mathbb{D}}$.

We shall denote by $\mathcal{C}^{\omega}$ the class of real-analytic functions. Further, saying that $f$ is of class $\mathcal{C}^{k}(\mathbb{T}), k \in(0, \infty] \cup\{\omega\}$, we mean that the function $t \mapsto$ $f\left(e^{i t}\right), t \in \mathbb{R}$, is in $\mathcal{C}^{k}(\mathbb{R})$. For a compact set $K \subset \mathbb{C}^{n}$ let $\mathcal{O}(K)$ denote the set of functions extending holomorphically on a neighborhood of $K$ (we assume that all neighborhoods are open). In that case we shall sometimes say that a given function is of class $\mathcal{O}(K)$. Note that $\mathcal{C}^{\omega}(\mathbb{T})=\mathcal{O}(\mathbb{T})$.

Let $|\cdot|$ denote the Euclidean norm in $\mathbb{C}^{n}$ and let $\operatorname{dist}(z, S):=\inf \{|z-s|$ : $s \in S\}$ be the distance of the point $z \in \mathbb{C}^{n}$ to a set $S \subset \mathbb{C}^{n}$. Moreover, for $S \subset \mathbb{C}^{n}$ we define $S_{*}:=S \backslash\{0\}$. Let $\mathbb{B}_{n}:=\left\{z \in \mathbb{C}^{n}:|z|=1\right\}$ be the unit ball and $B_{n}(a, r):=\left\{z \in \mathbb{C}^{n}:|z-a|<r\right\}$ the open ball with center $a \in \mathbb{C}^{n}$ and radius $r>0$. Put

$$
z \bullet w:=\sum_{j=1}^{n} z_{j} w_{j}
$$

for $z, w \in \mathbb{C}^{n}$ and let $\langle\cdot,-\rangle$ be the hermitian inner product on $\mathbb{C}^{n}$. The real inner product on $\mathbb{C}^{n}$ is denoted by $\langle\cdot,-\rangle_{\mathbb{R}}=\operatorname{Re}\langle\cdot,-\rangle$.

We use $\nabla$ to denote the gradient $\left(\partial / \partial x_{1}, \ldots, \partial / \partial x_{2 n}\right)$. For real-valued functions the gradient is naturally identified with $2\left(\partial / \partial \bar{z}_{1}, \ldots, \partial / \partial \bar{z}_{n}\right)$. Recall that

$$
\nu_{D}(a)=\frac{\nabla r(a)}{|\nabla r(a)|} .
$$

Let $\mathcal{H}$ be the Hessian matrix

$$
\left[\frac{\partial^{2}}{\partial x_{j} \partial x_{k}}\right]_{1 \leq j, k \leq 2 n}
$$

Sometimes, for a $\mathcal{C}^{2}$-smooth function $u$ and a vector $X \in \mathbb{R}^{2 n}$ the Hessian

$$
\sum_{j, k=1}^{2 n} \frac{\partial^{2} u}{\partial x_{j} \partial x_{k}}(a) X_{j} X_{k}=X^{T} \mathcal{H} u(a) X
$$

will be denoted by $\mathcal{H} u(a ; X)$. We denote by $\|\cdot\|$ the operator norm.
Definition 1.1. Let $D \subset \mathbb{C}^{n}$ be a domain. We say that $D$ is linearly convex (resp. weakly linearly convex) if through any point $a \in \mathbb{C}^{n} \backslash D$ (resp. $a \in \partial D)$ there passes an $(n-1)$-dimensional complex hyperplane disjoint from $D$.

A domain $D$ is said to be strongly linearly convex if

- $D$ has $\mathcal{C}^{2}$-smooth boundary;
- there exists a defining function $r$ of $D$ such that

$$
\begin{equation*}
\sum_{j, k=1}^{n} \frac{\partial^{2} r}{\partial z_{j} \partial \bar{z}_{k}}(a) X_{j} \bar{X}_{k}>\left|\sum_{j, k=1}^{n} \frac{\partial^{2} r}{\partial z_{j} \partial z_{k}}(a) X_{j} X_{k}\right|, \quad a \in \partial D, X \in T_{D}^{\mathbb{C}}(a)_{*} \tag{5}
\end{equation*}
$$

More generally, any point $a \in \partial D$ for which there exists a defining function $r$ satisfying (5) is called a point of strong linear convexity of $D$.

Furthermore, we say that a domain $D$ has real-analytic boundary if it possesses a real-analytic defining function.

Note that the condition (5) does not depend on the choice of the defining function of $D$.

REMARK 1.2 . Let $D \subset \mathbb{C}^{n}$ be a strongly linearly convex domain. Then

- any ( $n-1$ )-dimensional complex tangent hyperplane intersects $\partial D$ at precisely one point; in other words

$$
\bar{D} \cap\left(a+T_{D}^{\mathbb{C}}(a)\right)=\{a\}, \quad a \in \partial D
$$

- for $a \in \partial D$ the equation $\left\langle w-a, \nu_{D}(a)\right\rangle=0$ describes the $(n-1)$ dimensional complex tangent hyperplane $a+T_{D}^{\mathbb{C}}(a)$, consequently

$$
\left\langle z-a, \nu_{D}(a)\right\rangle \neq 0, \quad z \in D, a \in \partial D
$$

The main aim of the paper is to present a detailed proof of
Theorem 1.3 (Lempert Theorem). Let $D \subset \mathbb{C}^{n}$, $n \geq 2$, be a bounded strongly linearly convex domain. Then

$$
c_{D}=\widetilde{k}_{D} \quad \text { and } \quad \gamma_{D}=\kappa_{D}
$$

Important roles will be played by strongly convex domains and strongly convex functions.

Definition 1.4. A domain $D \subset \mathbb{C}^{n}$ is called strongly convex if

- $D$ has $\mathcal{C}^{2}$-smooth boundary;
- there exists a defining function $r$ of $D$ such that

$$
\begin{equation*}
\sum_{j, k=1}^{2 n} \frac{\partial^{2} r}{\partial x_{j} \partial x_{k}}(a) X_{j} X_{k}>0, \quad a \in \partial D, X \in T_{D}^{\mathbb{R}}(a)_{*} \tag{6}
\end{equation*}
$$

Generally, any point $a \in \partial D$ for which there exists a defining function $r$ satisfying (6) is called a point of strong convexity of $D$.

REMARK 1.5. A strongly convex domain $D \subset \mathbb{C}^{n}$ is convex and strongly linearly convex. Moreover, it is strictly convex, i.e. for any different points
$a, b \in \bar{D}$ the interior of the segment $[a, b]=\{t a+(1-t) b: t \in[0,1]\}$ is contained in $D$ (i.e. $t a+(1-t) b \in D$ for any $t \in(0,1)$ ).

Observe also that any bounded convex domain with real-analytic boundary is strictly convex. Actually, if a domain $D$ with real-analytic boundary were not strictly convex, then we would be able to find two distinct points $a, b \in \partial D$ such that the segment $[a, b]$ lies entirely in $\partial D$. On the other hand, the identity principle would imply that the set

$$
\{t \in \mathbb{R}: \exists \varepsilon>0: s a+(1-s) b \in \partial D \text { for }|s-t|<\varepsilon\}
$$

is open-closed in $\mathbb{R}$. Therefore it has to be empty. This immediately gives a contradiction.

REMARK 1.6. It is well-known that for any convex domain $D \subset \mathbb{C}^{n}$ there is a sequence $\left\{D_{m}\right\}$ of bounded strongly convex domains with real-analytic boundaries and such that $D_{m} \subset D_{m+1}$ and $\bigcup_{m} D_{m}=D$.

In particular, Theorem 1.3 holds for convex domains.
Definition 1.7. Let $U \subset \mathbb{C}^{n}$ be a domain. A function $u: U \rightarrow \mathbb{R}$ is called strongly convex if

- $u$ is $\mathcal{C}^{2}$-smooth;

$$
\sum_{j, k=1}^{2 n} \frac{\partial^{2} u}{\partial x_{j} \partial x_{k}}(a) X_{j} X_{k}>0, \quad a \in U, X \in\left(\mathbb{R}^{2 n}\right)_{*}
$$

Definition 1.8. The degree of a continuous function (treated as a curve) $\mathbb{T} \rightarrow \mathbb{T}$ is called its winding number. The fundamental group is a homotopy invariant. Thus the definition of the winding number of a continuous function $\varphi: \mathbb{T} \rightarrow \mathbb{C}_{*}$ is the same. We denote it by wind $\varphi$.

In the case of a $\mathcal{C}^{1}$-smooth function $\varphi: \mathbb{T} \rightarrow \mathbb{C}_{*}$, its winding number is just the index of $\varphi$ at 0 , i.e.

$$
\text { wind } \varphi=\frac{1}{2 \pi i} \int_{\varphi(\mathbb{T})} \frac{d \zeta}{\zeta}=\frac{1}{2 \pi i} \int_{0}^{2 \pi} \frac{\frac{d}{d t} \varphi\left(e^{i t}\right)}{\varphi\left(e^{i t}\right)} d t
$$

Remark 1.9.

- If $\varphi \in \mathcal{C}\left(\mathbb{T}, \mathbb{C}_{*}\right)$ extends to a function $\widetilde{\varphi} \in \mathcal{O}(\mathbb{D}) \cap \mathcal{C}(\overline{\mathbb{D}})$ then wind $\varphi$ is the number of zeroes of $\widetilde{\varphi}$ in $\mathbb{D}$ counted with multiplicities;
- $\operatorname{wind}(\varphi \psi)=$ wind $\varphi+\operatorname{wind} \psi$ for $\varphi, \psi \in \mathcal{C}\left(\mathbb{T}, \mathbb{C}_{*}\right)$;
- wind $\varphi=0$ if $\varphi \in \mathcal{C}(\mathbb{T})$ and $\operatorname{Re} \varphi>0$.

Definition 1.10. The boundary of a domain $D$ of $\mathbb{C}^{n}$ is real-analytic in a neighborhood $U$ of the set $S \subset \partial D$ if there exists a function $r \in \mathcal{C}^{\omega}(U, \mathbb{R})$ such that $D \cap U=\{z \in U: r(z)<0\}$ and $\nabla r$ does not vanish in $U$.

Definition 1.11. Let $D \subset \mathbb{C}^{n}$ be a domain with $\mathcal{C}^{1}$-smooth boundary. We call a holomorphic mapping $f: \mathbb{D} \rightarrow D$ a stationary mapping if
(1) $f$ extends to a holomorphic mapping (denoted by the same letter) in a neighborhood of $\overline{\mathbb{D}}$;
(2) $f(\mathbb{T}) \subset \partial D$;
(3) there exists a real-analytic function $\rho: \mathbb{T} \rightarrow \mathbb{R}_{>0}$ such that the mapping $\mathbb{T} \ni \zeta \mapsto \zeta \rho(\zeta) \overline{\nu_{D}(f(\zeta))} \in \mathbb{C}^{n}$ extends to a mapping (denoted by $\widetilde{f}$ ) holomorphic in a neighborhood of $\overline{\mathbb{D}}$.

Furthermore, we call a holomorphic mapping $f: \mathbb{D} \rightarrow D$ a weak stationary mapping if
$\left(1^{\prime}\right) f$ extends to a $\mathcal{C}^{1 / 2}{ }^{-}$-smooth mapping (denoted by the same letter) on $\overline{\mathbb{D}}$;
$\left(2^{\prime}\right) f(\mathbb{T}) \subset \partial D ;$
(3') there exists a $\mathcal{C}^{1 / 2}$-smooth function $\rho: \mathbb{T} \rightarrow \mathbb{R}_{>0}$ such that the mapping $\mathbb{T} \ni \zeta \mapsto \zeta \rho(\zeta) \overline{\nu_{D}(f(\zeta))} \in \mathbb{C}^{n}$ extends to a mapping $\widetilde{f} \in$ $\mathcal{O}(\mathbb{D}) \cap \mathcal{C}^{1 / 2}(\overline{\mathbb{D}})$.

The definition of a (weak) stationary mapping $f: \mathbb{D} \rightarrow D$ extends naturally to the case when $\partial D$ is real-analytic in a neighborhood of $f(\mathbb{T})$.

Directly from the definition of a stationary mapping $f$, it follows that $f$ and $\tilde{f}$ extend holomorphically on some neighborhoods of $\overline{\mathbb{D}}$. We shall denote their intersection by $\mathbb{D}_{f}$.

DEFINITION 1.12. Let $D \subset \mathbb{C}^{n}, n \geq 2$, be a bounded strongly linearly convex domain. A holomorphic mapping $f: \mathbb{D} \rightarrow D$ is called a (weak) E-mapping if it is a (weak) stationary mapping and
(4) setting $\varphi_{z}(\zeta):=\left\langle z-f(\zeta), \nu_{D}(f(\zeta))\right\rangle, \zeta \in \mathbb{T}$, we have wind $\varphi_{z}=0$ for some $z \in D$.

REmARK 1.13. The strong linear convexity of $D$ implies $\varphi_{z}(\zeta) \neq 0$ for any $z \in D$ and $\zeta \in \mathbb{T}$. Therefore, wind $\varphi_{z}$ vanishes for all $z \in D$ if it vanishes for some $z \in D$.

Additionally, any stationary mapping of a convex domain is an $E$-mapping $\left(\right.$ as $\left.\operatorname{Re} \varphi_{z}<0\right)$.

We shall prove that in a class of non-planar bounded strongly linearly convex domains with real-analytic boundaries, weak stationary mappings are just stationary mappings, so there is no difference between $E$-mappings and weak $E$-mappings.

We have the following result describing extremal mappings.

TheOrem 1.14. Let $D \subset \mathbb{C}^{n}$, $n \geq 2$, be a bounded strongly linearly convex domain. Then a holomorphic mapping $f: \mathbb{D} \rightarrow D$ is an extremal if and only if it is a weak E-mapping.

For a domain $D$ with real-analytic boundary, a holomorphic mapping $f: \mathbb{D} \rightarrow D$ is an extremal if and only if it is an E-mapping.

If $\partial D$ is of class $\mathcal{C}^{k}, k=3,4, \ldots, \infty$, then any weak $E$-mapping $f: \mathbb{D} \rightarrow D$ and its associated mappings $\widetilde{f}, \rho$ are $\mathcal{C}^{k-1-\varepsilon}{ }_{-s m o o t h ~ f o r ~ a n y ~} \varepsilon>0$.

The idea of the proof of the Lempert Theorem is as follows. In the realanalytic case we shall show that $E$-mappings are complex geodesics (because they have left inverses). Then we shall prove that for any different points $z, w \in D$ (resp. for a point $z \in D$ and a vector $v \in\left(\mathbb{C}^{n}\right)_{*}$ ) there is an $E$ mapping passing through $z, w$ (resp. such that $f(0)=z$ and $f^{\prime}(0)=v$ ). This will give the equality between the Lempert function and the Carathéodory distance. In the general case, we exhaust a $\mathcal{C}^{2}$-smooth domain by strongly linearly convex domains with real-analytic boundaries.

To prove Theorem 1.14 , we shall additionally observe that (weak) $E$ mappings are unique extremals.

In what follows and if not mentioned otherwise, $D \subset \mathbb{C}^{n}, n \geq 2$, is a bounded strongly linearly convex domain with real-analytic boundary.

## 2. Weak stationary mappings of strongly linearly convex domains with real-analytic boundaries are stationary mappings

Proposition 2.1. A weak stationary mapping of $D$ is a stationary mapping of $D$ with the same associated mappings.

Proof. Let $f: \mathbb{D} \rightarrow D$ be a weak stationary mapping. Our aim is to prove that $f, \widetilde{f} \in \mathcal{O}(\overline{\mathbb{D}})$ and $\rho \in \mathcal{C}^{\omega}(\mathbb{T})$.

Choose $\zeta_{0} \in \mathbb{T}$. Since $\widetilde{f}\left(\zeta_{0}\right) \neq 0$, we can assume that $\widetilde{f}_{1} \neq 0$ in $\overline{\mathbb{D}} \cap U_{0}$, where $U_{0}$ is a neighborhood of $\zeta_{0}$. This implies $\nu_{D, 1}\left(f\left(\zeta_{0}\right)\right) \neq 0$, so $\nu_{D, 1}$ does not vanish on some set $V_{0} \subset \partial D$, relatively open in $\partial D$, containing the point $f\left(\zeta_{0}\right)$. Shrinking $U_{0}$, if necessary, we may assume that $f\left(\mathbb{T} \cap U_{0}\right) \subset V_{0}$.

Define $\psi: V_{0} \rightarrow \mathbb{C}^{2 n-1}$ by

$$
\psi(z)=\left(z_{1}, \ldots, z_{n}, \overline{\left(\frac{\nu_{D, 2}(z)}{\nu_{D, 1}(z)}\right)}, \ldots, \overline{\left(\frac{\nu_{D, n}(z)}{\nu_{D, 1}(z)}\right)}\right) .
$$

The set $M:=\psi\left(V_{0}\right)$ is the graph of a $\mathcal{C}^{\omega}$ function defined on the local $\mathcal{C}^{\omega}$ submanifold $V_{0}$, so it is a local $\mathcal{C}^{\omega}$ submanifold in $\mathbb{C}^{2 n-1}$ of real dimension $2 n-1$. Direct calculations show that $M$ is totally real.

Let

$$
g(\zeta):=\left(f_{1}(\zeta), \ldots, f_{n}(\zeta), \frac{\widetilde{f}_{2}(\zeta)}{\widetilde{f}_{1}(\zeta)}, \ldots, \frac{\widetilde{f}_{n}(\zeta)}{\widetilde{f}_{1}(\zeta)}\right), \quad \zeta \in \overline{\mathbb{D}} \cap U_{0}
$$

If $\zeta \in \mathbb{T} \cap U_{0}$ then $\widetilde{f}_{k}(\zeta) \widetilde{f}_{1}(\zeta)^{-1}=\overline{\nu_{D, k}(f(\zeta))} \overline{\nu_{D, 1}(f(\zeta))}{ }^{-1}$, so $g(\zeta)=\psi(f(\zeta))$. Therefore, $g\left(\mathbb{T} \cap U_{0}\right) \subset M$. Thanks to the Reflection Principle (see Appendix), $g$ extends holomorphically past $\mathbb{T} \cap U_{0}$, so $f$ extends holomorphically on a neighborhood of $\zeta_{0}$.

The mapping $\overline{\nu_{D} \circ f}$ is real-analytic on $\mathbb{T}$, so it extends to a mapping $h$ holomorphic in a neighborhood $W$ of $\mathbb{T}$. For $\zeta \in \mathbb{T} \cap U_{0}$ we have

$$
\frac{\zeta h_{1}(\zeta)}{\widetilde{f}_{1}(\zeta)}=\frac{1}{\rho(\zeta)}
$$

The function on the left side is holomorphic in $\mathbb{D} \cap U_{0} \cap W$ and continuous in $\overline{\mathbb{D}} \cap U_{0} \cap W$. Since it has real values on $\mathbb{T} \cap U_{0}$, the Reflection Principle implies that it is holomorphic in a neighborhood of $\mathbb{T} \cap U_{0}$. Hence $\rho$ and $\widetilde{f}$ are holomorphic in a neighborhood of $\zeta_{0}$.

## 3. (Weak) E-mappings vs. extremal mappings and complex geo-

 desics. In this section we will prove some important properties of (weak) $E$-mappings. In particular, we will show that they are complex geodesics and unique extremals.3.1. Weak $E$-mappings are complex geodesics and unique extremals. The results of this subsection concern weak $E$-mappings of bounded strongly linearly convex domains $D \subset \mathbb{C}^{n}, n \geq 2$.

Let

$$
G(z, \zeta):=(z-f(\zeta)) \bullet \tilde{f}(\zeta), \quad z \in \mathbb{C}^{n}, \zeta \in \mathbb{D}_{f}
$$

Proposition 3.1.1. Let $D \subset \mathbb{C}^{n}, n \geq 2$, be a bounded strongly linearly convex domain and let $f: \mathbb{D} \rightarrow D$ be a weak E-mapping. Then there exist an open set $W \supset \bar{D} \backslash f(\mathbb{T})$ and a holomorphic mapping $F: W \rightarrow \mathbb{D}$ such that for any $z \in W$ the number $F(z)$ is a unique solution of the equation $G(z, \zeta)=0, \zeta \in \mathbb{D}$. In particular, $F \circ f=\mathrm{id}_{\mathbb{D}}$.

In what follows we will strengthen the above proposition for domains with real-analytic boundaries (see Proposition 3.2.4).

Proof of Proposition 3.1.1. Set $A:=\bar{D} \backslash f(\mathbb{T})$. Since $D$ is strongly linearly convex, $\varphi_{z}$ does not vanish in $\mathbb{T}$ for any $z \in A$, so by a continuity argument the condition (4) of Definition 1.12 holds for every $z$ in some open set $W \supset A$. For a fixed $z \in W$ we have

$$
G(z, \zeta)=\zeta \rho(\zeta) \varphi_{z}(\zeta), \quad \zeta \in \mathbb{T},
$$

so wind $G(z, \cdot)=1$. Since $G(z, \cdot) \in \mathcal{O}(\mathbb{D})$, it has in $\mathbb{D}$ exactly one simple root $F(z)$. Hence $G(z, F(z))=0$ and $\frac{\partial G}{\partial \zeta}(z, F(z)) \neq 0$. By the Implicit Function Theorem, $F$ is holomorphic in $W$. The equality $F(f(\zeta))=\zeta$ for $\zeta \in \mathbb{D}$ is clear.

From the proposition above we immediately get
Corollary 3.1.2. Every weak E-mapping $f: \mathbb{D} \rightarrow D$ of a bounded strongly linearly convex domain $D \subset \mathbb{C}^{n}, n \geq 2$, is a complex geodesic. In particular,

$$
c_{D}(f(\zeta), f(\xi))=\widetilde{k}_{D}(f(\zeta), f(\xi)) \quad \text { and } \quad \gamma_{D}\left(f(\zeta) ; f^{\prime}(\zeta)\right)=\kappa_{D}\left(f(\zeta) ; f^{\prime}(\zeta)\right)
$$

for any $\zeta, \xi \in \mathbb{D}$.
Using left inverses of weak $E$-mappings we may prove the uniqueness of extremals.

Proposition 3.1.3. Let $D \subset \mathbb{C}^{n}, n \geq 2$, be a bounded strongly linearly convex domain and let $f: \mathbb{D} \rightarrow D$ be a weak $E$-mapping. Then for any $\xi \in(0,1)$ the mapping $f$ is a unique $\widetilde{k}_{D}$-extremal for $z=f(0), w=f(\xi)$ (resp. a unique $\kappa_{D}$-extremal for $z=f(0), v=f^{\prime}(0)$ ).

Proof. Suppose that $g$ is a $\widetilde{k}_{D}$-extremal for $z, w$ (resp. a $\kappa_{D}$-extremal for $z, v)$ such that $g(0)=z, g(\xi)=w$ (resp. $g(0)=z, g^{\prime}(0)=v$ ). Our aim is to show that $f=g$. Proposition 3.1.1 provides us with the mapping $F$, which is a left inverse for $f$. By the Schwarz Lemma, $F$ is a left inverse for $g$ as well, that is, $F \circ g=\mathrm{id}_{\mathbb{D}}$. We claim that $\lim _{\mathbb{D} \ni \zeta \rightarrow \zeta_{0}} g(\zeta)=f\left(\zeta_{0}\right)$ for any $\zeta_{0} \in \mathbb{T}$ (in particular, we shall show that the limit exists).

Assume the contrary. Then there are $\zeta_{0} \in \mathbb{T}$ and a sequence $\left\{\zeta_{m}\right\} \subset \mathbb{D}$ convergent to $\zeta_{0}$ such that the limit $Z:=\lim _{m \rightarrow \infty} g\left(\zeta_{m}\right) \in \bar{D}$ exists and is not equal to $f\left(\zeta_{0}\right)$. We have $G(z, F(z))=0$, so putting $z=g\left(\zeta_{m}\right)$ we infer that

$$
0=\left(g\left(\zeta_{m}\right)-f\left(F\left(g\left(\zeta_{m}\right)\right)\right)\right) \bullet \widetilde{f}\left(F\left(g\left(\zeta_{m}\right)\right)\right)=\left(g\left(\zeta_{m}\right)-f\left(\zeta_{m}\right)\right) \bullet \tilde{f}\left(\zeta_{m}\right)
$$

Letting $m \rightarrow \infty$ we get

$$
0=\left(Z-f\left(\zeta_{0}\right)\right) \bullet \tilde{f}\left(\zeta_{0}\right)=\zeta_{0} \rho\left(\zeta_{0}\right)\left\langle Z-f\left(\zeta_{0}\right), \nu_{D}\left(f\left(\zeta_{0}\right)\right)\right\rangle
$$

This means that $Z-f\left(\zeta_{0}\right) \in T_{D}^{\mathbb{C}}\left(f\left(\zeta_{0}\right)\right)$. Since $D$ is strongly linearly convex, we deduce that $Z=f\left(\zeta_{0}\right)$, which is a contradiction.

Hence $g$ extends continuously on $\overline{\mathbb{D}}$ and, by the maximum principle, $g=f$.

Proposition 3.1.4. Let $D \subset \mathbb{C}^{n}, n \geq 2$, be a bounded strongly linearly convex domain, let $f: \mathbb{D} \rightarrow D$ be a weak $E$-mapping and let $a$ be an automorphism of $\mathbb{D}$. Then $f \circ a$ is a weak $E$-mapping of $D$.

Proof. Set $g:=f \circ a$.
Clearly, the conditions ( $1^{\prime}$ ) and ( $2^{\prime}$ ) of Definition 1.11 are satisfied by $g$.
To prove that $g$ satisfies the condition (4) of Definition 1.12 fix a point $z \in D$. Let $\varphi_{z, f}, \varphi_{z, g}$ be the functions appearing in the condition (4) for $f, g$ respectively. Then $\varphi_{z, g}=\varphi_{z, f} \circ a$. Since $a$ maps $\mathbb{T}$ onto $\mathbb{T}$ diffeomorphically, we have wind $\varphi_{z, g}= \pm$ wind $\varphi_{z, f}=0$.

It remains to show that the condition $\left(3^{\prime}\right)$ of Definition 1.11 is also satisfied by $g$. Note that the function $\widetilde{a}(\zeta):=\zeta / a(\zeta)$ has a holomorphic branch of the logarithm in the neighborhood of $\mathbb{T}$. This follows from wind $\widetilde{a}=0$, but it may also be shown quite elementarily. Actually, it suffices to prove that $\widetilde{a}(\mathbb{T}) \neq \mathbb{T}$. Expand $a$ as

$$
a(\zeta)=e^{i t} \frac{\zeta-b}{1-\bar{b} \zeta}
$$

with some $t \in \mathbb{R}, b \in \mathbb{D}$, and observe that $\widetilde{a}$ does not attain the value $-e^{-i t}$. Indeed, if $\zeta / a(\zeta)=-e^{-i t}$ for some $\zeta \in \mathbb{T}$, then

$$
\frac{1-\bar{b} \zeta}{1-b \bar{\zeta}}=-1
$$

so $2=2 \operatorname{Re}(b \bar{\zeta}) \leq 2|b|$, which is impossible.
Concluding, there exists a function $v$ holomorphic in a neighborhood of $\mathbb{T}$ such that

$$
\frac{\zeta}{a(\zeta)}=e^{i v(\zeta)}
$$

Note that $v(\mathbb{T}) \subset \mathbb{R}$. Expanding $v$ in Laurent series

$$
v(\zeta)=\sum_{k=-\infty}^{\infty} a_{k} \zeta^{k}, \quad \zeta \text { near } \mathbb{T},
$$

we infer that $a_{-k}=\bar{a}_{k}, k \in \mathbb{Z}$. Therefore,

$$
v(\zeta)=a_{0}+\sum_{k=1}^{\infty} 2 \operatorname{Re}\left(a_{k} \zeta^{k}\right)=\operatorname{Re}\left(a_{0}+2 \sum_{k=1}^{\infty} a_{k} \zeta^{k}\right), \quad \zeta \in \mathbb{T} .
$$

Hence, there is a function $h$ holomorphic in a neighborhood of $\overline{\mathbb{D}}$ such that $v=\operatorname{Im} h$. Put $u:=h-i v$. Then $u \in \mathcal{O}(\mathbb{T})$ and $u(\mathbb{T}) \subset \mathbb{R}$.

Take $\rho$ as in the condition ( $3^{\prime}$ ) of Definition 1.11 for $f$ and define

$$
r(\zeta):=\rho(a(\zeta)) e^{u(\zeta)}, \quad \zeta \in \mathbb{T}
$$

We compute

$$
\begin{aligned}
\zeta r(\zeta) \overline{\nu_{D}(g(\zeta))} & =\zeta u^{u(\zeta)} \rho(a(\zeta)) \overline{\nu_{D}(f(a(\zeta)))} \\
& =a(\zeta) h(\zeta) \rho(a(\zeta)) \overline{\nu_{D}(f(a(\zeta)))}=h(\zeta) \tilde{f}(a(\zeta)), \quad \zeta \in \mathbb{T} .
\end{aligned}
$$

Thus $\mathbb{T} \ni \zeta \mapsto \zeta r(\zeta) \overline{\nu_{D}(g(\zeta))} \in \mathbb{C}^{n}$ extends to a mapping from $\mathcal{O}(\mathbb{D}) \cap$ $\mathcal{C}^{1 / 2}(\overline{\mathbb{D}})$.

Corollary 3.1.5. Every weak E-mapping $f: \mathbb{D} \rightarrow D$ of a bounded strongly linearly convex domain $D \subset \mathbb{C}^{n}, n \geq 2$, is a unique $\widetilde{k}_{D}$-extremal for $f(\zeta), f(\xi)$ (resp. a unique $\kappa_{D}$-extremal for $f(\zeta), f^{\prime}(\zeta)$ ), for any $\zeta, \xi \in \mathbb{D}$, $\zeta \neq \xi$.
3.2. Generalization of Proposition 3.1.1. The results obtained in this subsection will play an important role later.

We start with
Proposition 3.2.1. Let $f: \mathbb{D} \rightarrow D$ be an E-mapping. Then the function $f^{\prime} \bullet \tilde{f}$ is a positive constant.

Proof. Consider the curve

$$
\mathbb{R} \ni t \mapsto f\left(e^{i t}\right) \in \partial D
$$

The tangent vector $i e^{i t} f^{\prime}\left(e^{i t}\right)$ to this curve belongs to $T_{D}^{\mathbb{R}}\left(f\left(e^{i t}\right)\right)$, i.e.

$$
\operatorname{Re}\left\langle i e^{i t} f^{\prime}\left(e^{i t}\right), \nu_{D}\left(f\left(e^{i t}\right)\right)\right\rangle=0
$$

Thus for $\zeta \in \mathbb{T}$,

$$
0=\rho(\zeta) \operatorname{Re}\left\langle i \zeta f^{\prime}(\zeta), \nu_{D}(f(\zeta))\right\rangle=-\operatorname{Im} f^{\prime}(\zeta) \bullet \widetilde{f}(\zeta)
$$

so the holomorphic function $f^{\prime} \bullet \tilde{f}$ is a real constant $C$.
Considering the curve

$$
[0,1+\varepsilon) \ni t \mapsto f(t) \in \mathbb{C}^{n}
$$

for small $\varepsilon>0$ and noting that $f([0,1)) \subset D, f(1) \in \partial D$, we see that the derivative of $r \circ f$ at a point $t=1$ is non-negative, where $r$ is a defining function of $D$. Hence

$$
0 \leq \operatorname{Re}\left\langle f^{\prime}(1), \nu_{D}(f(1))\right\rangle=\frac{1}{\rho(1)} \operatorname{Re}\left(f^{\prime}(1) \bullet \widetilde{f}(1)\right)=\frac{C}{\rho(1)},
$$

i.e. $C \geq 0$. For $\zeta \in \mathbb{T}$,

$$
\frac{f(\zeta)-f(0)}{\zeta} \bullet \tilde{f}(\zeta)=\rho(\zeta)\left\langle f(\zeta)-f(0), \nu_{D}(f(\zeta))\right\rangle
$$

This function has winding number 0 . Therefore, the function

$$
g(\zeta):=\frac{f(\zeta)-f(0)}{\zeta} \bullet \widetilde{f}(\zeta)
$$

which is holomorphic in a neighborhood of $\overline{\mathbb{D}}$, does not vanish in $\mathbb{D}$. In particular, $C=g(0) \neq 0$.

Since the function $\rho$ is defined up to a constant factor, we can choose it so that $f^{\prime} \bullet \tilde{f} \equiv 1$, i.e.

$$
\begin{equation*}
\rho(\zeta)^{-1}=\left\langle\zeta f^{\prime}(\zeta), \nu_{D}(f(\zeta))\right\rangle, \quad \zeta \in \mathbb{T} \tag{7}
\end{equation*}
$$

In that way $\tilde{f}$ and $\rho$ are uniquely determined by $f$.
Proposition 3.2.2. Every E-mapping $f: \mathbb{D} \rightarrow D$ is injective in $\overline{\mathbb{D}}$.
Proof. The function $f$ has a left inverse in $\mathbb{D}$, so it suffices to check the injectivity on $\mathbb{T}$. Suppose that $f\left(\zeta_{1}\right)=f\left(\zeta_{2}\right)$ for some $\zeta_{1}, \zeta_{2} \in \mathbb{T}, \zeta_{1} \neq \zeta_{2}$,
and consider the curves

$$
\gamma_{j}:[0,1] \ni t \mapsto f\left(t \zeta_{j}\right) \in \bar{D}, \quad j=1,2
$$

Since

$$
\operatorname{Re}\left\langle\gamma_{j}^{\prime}(1), \nu_{D}\left(f\left(\zeta_{j}\right)\right)\right\rangle=\operatorname{Re}\left\langle\zeta_{j} f^{\prime}\left(\zeta_{j}\right), \nu_{D}\left(f\left(\zeta_{j}\right)\right)\right\rangle=\rho\left(\zeta_{j}\right)^{-1} \neq 0
$$

the curves $\gamma_{j}$ hit $\partial D$ transversally at their common point $f\left(\zeta_{1}\right)$. We claim that there exists $C>0$ such that for $t \in(0,1)$ close to 1 there is $s_{t} \in(0,1)$ satisfying $\widetilde{k}_{D}\left(f\left(t \zeta_{1}\right), f\left(s_{t} \zeta_{2}\right)\right)<C$. This will finish the proof since

$$
\widetilde{k}_{D}\left(f\left(t \zeta_{1}\right), f\left(s_{t} \zeta_{2}\right)\right)=p\left(t \zeta_{1}, s_{t} \zeta_{2}\right) \rightarrow \infty, \quad t \rightarrow 1
$$

We may assume that $f\left(\zeta_{1}\right)=0$ and $\nu_{D}(0)=(1,0, \ldots, 0)=: e_{1}$. There exists a ball $B \subset D$ tangent to $\partial D$ at 0 . Using a homothety if necessary, one can assume that $B=\mathbb{B}_{n}-e_{1}$. From the transversality of $\gamma_{1}, \gamma_{2}$ to $\partial D$ there exists a cone

$$
A:=\left\{z \in \mathbb{C}^{n}:-\operatorname{Re} z_{1}>k|z|\right\}, \quad k>0
$$

such that $\gamma_{1}(t), \gamma_{2}(t) \in A \cap B$ if $t \in(0,1)$ is sufficiently close to 1 . For $z \in A$ let $k_{z}>k$ be the positive number satisfying

$$
|z|=\frac{-\operatorname{Re} z_{1}}{k_{z}}
$$

Note that for any $a \in \gamma_{1}((0,1))$ sufficiently close to 0 one may find $b \in$ $\gamma_{2}((0,1)) \cap A \cap B$ such that $\operatorname{Re} b_{1}=\operatorname{Re} a_{1}$. To get a contradiction it suffices to show that $\widetilde{k}_{D}(a, b)$ is bounded from above by a constant independent of $a$ and $b$.

We have the estimate

$$
\begin{aligned}
\widetilde{k}_{D}(a, b) & \leq \widetilde{k}_{\mathbb{B}_{n}-e_{1}}(a, b)=\widetilde{k}_{\mathbb{B}_{n}}\left(a+e_{1}, b+e_{1}\right) \\
& =\tanh ^{-1} \sqrt{1-\frac{\left(1-\left|a+e_{1}\right|^{2}\right)\left(1-\left|b+e_{1}\right|^{2}\right)}{\left|1-\left\langle a+e_{1}, b+e_{1}\right\rangle\right|^{2}}}
\end{aligned}
$$

The last expression is bounded from above if and only if

$$
\frac{\left(1-\left|a+e_{1}\right|^{2}\right)\left(1-\left|b+e_{1}\right|^{2}\right)}{\left|1-\left\langle a+e_{1}, b+e_{1}\right\rangle\right|^{2}}
$$

is bounded from below by some positive constant. We estimate

$$
\begin{aligned}
\frac{\left(1-\left|a+e_{1}\right|^{2}\right)\left(1-\left|b+e_{1}\right|^{2}\right)}{\left|1-\left\langle a+e_{1}, b+e_{1}\right\rangle\right|^{2}} & =\frac{\left(2 \operatorname{Re} a_{1}+|a|^{2}\right)\left(2 \operatorname{Re} b_{1}+|b|^{2}\right)}{\left|\langle a, b\rangle+a_{1}+\bar{b}_{1}\right|^{2}} \\
& =\frac{\left(2 \operatorname{Re} a_{1}+\frac{\left(\operatorname{Re} a_{1}\right)^{2}}{k_{a}^{2}}\right)\left(2 \operatorname{Re} a_{1}+\frac{\left(\operatorname{Re} a_{1}\right)^{2}}{k_{b}^{2}}\right)}{\left|\langle a, b\rangle+2 \operatorname{Re} a_{1}+i \operatorname{Im} a_{1}-i \operatorname{Im} b_{1}\right|^{2}} \\
& \geq \frac{\left(\operatorname{Re} a_{1}\right)^{2}\left(2+\frac{\operatorname{Re} a_{1}}{k_{a}^{2}}\right)\left(2+\frac{\operatorname{Re} a_{1}}{k_{b}^{2}}\right)}{2\left|\langle a, b\rangle+i \operatorname{Im} a_{1}-i \operatorname{Im} b_{1}\right|^{2}+2\left|2 \operatorname{Re} a_{1}\right|^{2}}
\end{aligned}
$$

$$
\begin{aligned}
& \geq \frac{\left(\operatorname{Re} a_{1}\right)^{2}\left(2+\frac{\operatorname{Re} a_{1}}{k_{a}^{2}}\right)\left(2+\frac{\operatorname{Re} a_{1}}{k_{b}^{2}}\right)}{2(|a||b|+|a|+|b|)^{2}+8\left(\operatorname{Re} a_{1}\right)^{2}} \\
& =\frac{\left(\operatorname{Re} a_{1}\right)^{2}\left(2+\frac{\operatorname{Re} a_{1}}{k_{a}^{2}}\right)\left(2+\frac{\operatorname{Re} a_{1}}{k_{b}^{2}}\right)}{2\left(\frac{\left(-\operatorname{Re} a_{1}\right)^{2}}{k_{a}^{2} k_{b}^{2}}-\frac{\operatorname{Re} a_{1}}{k_{a}}-\frac{\operatorname{Re} a_{1}}{k_{b}}\right)^{2}+8\left(\operatorname{Re} a_{1}\right)^{2}} \\
& =\frac{\left(2+\frac{\operatorname{Re} a_{1}}{k_{a}^{2}}\right)\left(2+\frac{\operatorname{Re} a_{1}}{k_{b}^{2}}\right)}{2\left(\frac{-\operatorname{Re} a_{1}}{k_{a}^{2} k_{b}^{2}}+\frac{1}{k_{a}}+\frac{1}{k_{b}}\right)^{2}+8}>\frac{1}{2(1+2 / k)^{2}+8} .
\end{aligned}
$$

Assume that we are in the setting of Proposition 3.1.1 and $D$ has realanalytic boundary. Our aim is to replace $W$ with a neighborhood of $\bar{D}$.

REMARK 3.2.3. For $\zeta_{0} \in \mathbb{D}_{f}$ we have $G\left(f\left(\zeta_{0}\right), \zeta_{0}\right)=0$ and $\frac{\partial G}{\partial \zeta}\left(f\left(\zeta_{0}\right), \zeta_{0}\right)$ $=-1$. By the Implicit Function Theorem there exist neighborhoods $U_{\zeta_{0}}, V_{\zeta_{0}}$ of $f\left(\zeta_{0}\right), \zeta_{0}$ respectively and a holomorphic function $F_{\zeta_{0}}: U_{\zeta_{0}} \rightarrow V_{\zeta_{0}}$ such that for any $z \in U_{\zeta_{0}}$ the point $F_{\zeta_{0}}(\zeta)$ is the unique solution of the equation $G(z, \zeta)=0, \zeta \in V_{\zeta_{0}}$.

In particular, if $\zeta_{0} \in \mathbb{D}$ then $F_{\zeta_{0}}=F$ near $f\left(\zeta_{0}\right)$.
Proposition 3.2.4. Let $f: \mathbb{D} \rightarrow D$ be an E-mapping. Then there exist arbitrarily small neighborhoods $U, V$ of $\bar{D}, \overline{\mathbb{D}}$ respectively such that for any $z \in U$ the equation $G(z, \zeta)=0, \zeta \in V$, has exactly one solution.

Proof. In view of Proposition 3.1.1 and Remark 3.2.3, it suffices to prove that there exist neighborhoods $U, V$ of $\bar{D}, \overline{\mathbb{D}}$ respectively such that for any $z \in U$ the equation $G(z, \cdot)=0$ has at most one solution $\zeta \in V$.

Assume the contrary. Then for any neighborhoods $U$ of $\bar{D}$ and $V$ of $\overline{\mathbb{D}}$ there are $z \in U, \zeta_{1}, \zeta_{2} \in V, \zeta_{1} \neq \zeta_{2}$ such that $G\left(z, \zeta_{1}\right)=G\left(z, \zeta_{2}\right)=0$. For $m \in \mathbb{N}$ put

$$
U_{m}:=\left\{z \in \mathbb{C}^{n}: \operatorname{dist}(z, D)<1 / m\right\}, \quad V_{m}:=\{\zeta \in \mathbb{C}: \operatorname{dist}(\zeta, \mathbb{D})<1 / m\} .
$$

There exist $z_{m} \in U_{m}, \zeta_{m, 1}, \zeta_{m, 2} \in V_{m}, \zeta_{m, 1} \neq \zeta_{m, 2}$ such that $G\left(z_{m}, \zeta_{m, 1}\right)=$ $G\left(z_{m}, \zeta_{m, 2}\right)=0$. Passing to a subsequence we may assume that $z_{m} \rightarrow z_{0}$ $\in \bar{D}$. Analogously we may assume $\zeta_{m, 1} \rightarrow \zeta_{1} \in \overline{\mathbb{D}}$ and $\zeta_{m, 2} \rightarrow \zeta_{2} \in \overline{\mathbb{D}}$. Clearly, $G\left(z_{0}, \zeta_{1}\right)=G\left(z_{0}, \zeta_{2}\right)=0$. Let us consider the possible cases.

1) If $\zeta_{1}, \zeta_{2} \in \mathbb{T}$ then $G\left(z_{0}, \zeta_{j}\right)=0$ is equivalent to

$$
\left\langle z_{0}-f\left(\zeta_{j}\right), \nu_{D}\left(f\left(\zeta_{j}\right)\right)\right\rangle=0, \quad j=1,2
$$

consequently $z_{0}-f\left(\zeta_{j}\right) \in T_{D}^{\mathbb{C}}\left(f\left(\zeta_{j}\right)\right)$. By the strong linear convexity of $D$ we get $z_{0}=f\left(\zeta_{j}\right)$. But $f$ is injective in $\overline{\mathbb{D}}$, so $\zeta_{1}=\zeta_{2}=$ : $\zeta_{0}$. It follows from Remark 3.2 .3 that in a sufficiently small neighborhood of $\left(z_{0}, \zeta_{0}\right)$ all solutions of the equation $G(z, \zeta)=0$ are of the form $\left(z, F_{\zeta_{0}}(z)\right)$. The points $\left(z_{m}, \zeta_{m, 1}\right)$ and $\left(z_{m}, \zeta_{m, 2}\right)$ belong to this neighborhood for large $m$, which gives a contradiction.
2) If $\zeta_{1} \in \mathbb{T}$ and $\zeta_{2} \in \mathbb{D}$, then analogously to the above we deduce that $z_{0}=f\left(\zeta_{1}\right)$. Let us take an arbitrary sequence $\left\{\eta_{m}\right\} \subset \mathbb{D}$ convergent to $\zeta_{1}$. Then $f\left(\eta_{m}\right) \in D$ and $f\left(\eta_{m}\right) \rightarrow z_{0}$, so the sequence $G\left(f\left(\eta_{m}\right), \cdot\right)$ converges to $G\left(z_{0}, \cdot\right)$ uniformly on $\mathbb{D}$. Since $G\left(z_{0}, \cdot\right) \not \equiv 0, G\left(z_{0}, \zeta_{2}\right)=0$ and $\zeta_{2} \in \mathbb{D}$, we deduce from the Hurwitz Theorem that for large $m$ the functions $G\left(f\left(\eta_{m}\right), \cdot\right)$ have roots $\theta_{m} \in \mathbb{D}$ such that $\theta_{m} \rightarrow \zeta_{2}$. Hence $G\left(f\left(\eta_{m}\right), \theta_{m}\right)=0$ and from the uniqueness of solutions in $D \times \mathbb{D}$ (Proposition 3.1.1) we have

$$
\theta_{m}=F\left(f\left(\eta_{m}\right)\right)=\eta_{m}
$$

This is a contradiction, because the left side tends to $\zeta_{2}$ and the right side to $\zeta_{1}$, as $m \rightarrow \infty$.
3) We are left with the case $\zeta_{1}, \zeta_{2} \in \mathbb{D}$. If $z_{0} \in \bar{D} \backslash f(\mathbb{T})$ then $z_{0} \in W$. In $W \times \mathbb{D}$ all solutions of the equation $G=0$ are of the form $(z, F(z)), z \in W$. But for large $m$ the points $\left(z_{m}, \zeta_{m, 1}\right),\left(z_{m}, \zeta_{m, 2}\right)$ belong to $W \times \mathbb{D}$, which is a contradiction with the uniqueness.

If $z_{0} \in f(\mathbb{T})$, then $z_{0}=f\left(\zeta_{0}\right)$ for some $\zeta_{0} \in \mathbb{T}$. Clearly, $G\left(f\left(\zeta_{0}\right), \zeta_{0}\right)=0$, whence $G\left(z_{0}, \zeta_{0}\right)=G\left(z_{0}, \zeta_{1}\right)=0$ and $\zeta_{0} \in \mathbb{T}, \zeta_{1} \in \mathbb{D}$. This is just case 2$)$, considered above.

Corollary 3.2.5. There are neighborhoods $U, V$ of $\bar{D}, \overline{\mathbb{D}}$ respectively with $V \Subset \mathbb{D}_{f}$ such that the function $F$ extends holomorphically on $U$ and all solutions of the equation $\left.G\right|_{U \times V}=0$ are of the form $(z, F(z)), z \in U$. In particular, $F \circ f=\mathrm{id}_{V}$.

## 4. Hölder estimates

Definition 4.1. For a given $c>0$ let the family $\mathcal{D}(c)$ consist of all pairs $(D, z)$, where $D \subset \mathbb{C}^{n}, n \geq 2$, is a bounded strongly pseudoconvex domain with real-analytic boundary and $z \in D$, satisfying
(1) $\operatorname{dist}(z, \partial D) \geq 1 / c$;
(2) the diameter of $D$ is not greater than $c$ and $D$ satisfies the interior ball condition with a radius $1 / c$;
(3) for any $x, y \in D$ there exist $m \leq 8 c^{2}$ and open balls $B_{0}, \ldots, B_{m} \subset D$ of radius $1 /(2 c)$ such that $x \in B_{0}, y \in B_{m}$ and the distance between the centers of the balls $B_{j}, B_{j+1}$ is not greater than $1 /(4 c)$ for $j=$ $0, \ldots, m-1$;
(4) for any open ball $B \subset \mathbb{C}^{n}$ of radius not greater than $1 / c$, intersecting non-emptily with $\partial D$, there exists a mapping $\Phi \in \mathcal{O}\left(\bar{D}, \mathbb{C}^{n}\right)$ such that
(a) for any $w \in \Phi(B \cap \partial D)$ there is a ball of a radius containing $\Phi(D)$ and tangent to $\partial \Phi(D)$ at $w$ (let us call this the "exterior ball condition" with radius $c$ );
(b) $\Phi$ is biholomorphic in a neighborhood of $\bar{B}$ and $\Phi^{-1}(\Phi(B))=B$;
(c) the entries of all matrices $\Phi^{\prime}$ on $B \cap \bar{D}$ and $\left(\Phi^{-1}\right)^{\prime}$ on $\Phi(B \cap \bar{D})$ are bounded in modulus by $c$;
(d) $\operatorname{dist}(\Phi(z), \partial \Phi(D)) \geq 1 / c$;
(5) the normal vector $\nu_{D}$ is Lipschitz with constant $2 c$, that is,

$$
\left|\nu_{D}(a)-\nu_{D}(b)\right| \leq 2 c|a-b|, \quad a, b \in \partial D ;
$$

(6) the $\varepsilon$-hull of $D$, i.e. the domain $D_{\varepsilon}:=\left\{w \in \mathbb{C}^{n}: \operatorname{dist}(w, D)<\varepsilon\right\}$, is strongly pseudoconvex for any $\varepsilon \in(0,1 / c]$.

Recall that the interior ball condition with radius $r>0$ means that for any point $a \in \partial D$ there is $a^{\prime} \in D$ and a ball $B_{n}\left(a^{\prime}, r\right) \subset D$ tangent to $\partial D$ at $a$. Equivalently

$$
D=\bigcup_{a^{\prime} \in D^{\prime}} B_{n}\left(a^{\prime}, r\right)
$$

for some set $D^{\prime} \subset D$.
It may be shown that (2) and (5) may be expressed in terms of boundedness of the normal curvature, boundedness of the domain and the condition (3). This however lies beyond the scope of this paper and needs some very technical arguments so we omit the proof of this fact. The reason to use (2) in such a form is its connection with (3) (this allows us to simplify the proof in some points).

Remark 4.2. Note that any convex domain satisfying conditions (1)-(4) of Definition 4.1 satisfies conditions (5) and (6), as well.

Actually, it follows from (2) that for any $a \in \partial D$ there exists a ball $B_{n}\left(a^{\prime}, 1 / c\right) \subset D$ tangent to $\partial D$ at $a$. Then

$$
\nu_{D}(a)=\frac{a^{\prime}-a}{\left|a^{\prime}-a\right|}=c\left(a^{\prime}-a\right) .
$$

Hence
$\left|\nu_{D}(a)-\nu_{D}(b)\right|=c\left|a^{\prime}-a-b^{\prime}+b\right|=c\left|a^{\prime}-b^{\prime}-(a-b)\right| \leq c\left|a^{\prime}-b^{\prime}\right|+c|a-b|$.
Since $D$ is convex, we have $\left|a^{\prime}-b^{\prime}\right| \leq|a-b|$, which gives (5).
The condition (6) is also clear-for any $\varepsilon>0$ the $\varepsilon$-hull of a strongly convex domain is strongly convex.

Remark 4.3. For a convex domain $D$ the condition (3) of Definition 4.1 amounts to the condition (2).

Indeed, for two points $x, y \in D$ take two balls of radius $1 /(2 c)$ containing $x, y$ and contained in $D$. Then divide the interval between the centers of the balls into $\left[4 c^{2}\right]+1$ equal parts and take balls of radius $1 /(2 c)$ with centers at the points of the partition.

Note also that if a strongly convex domain $D$ satisfies the interior ball condition with radius $1 / c$ and the exterior ball condition with radius $c$, then one can take $\Phi:=\mathrm{id}_{\mathbb{C}^{n}}$.

REMARK 4.4. For a strongly pseudoconvex domain $D$ and $c^{\prime}>0$ and for any $z \in D$ such that $\operatorname{dist}(z, \partial D)>1 / c^{\prime}$ there exists $c=c\left(c^{\prime}\right)>0$ satisfying $(D, z) \in \mathcal{D}(c)$.

Indeed, the conditions (1)-(3) and (5)-(6) are clear. Only (4) is nontrivial. The construction of the mapping $\Phi$ amounts to constructing Fornæss peak functions. Actually, we can apply directly Proposition 1 from [2] to any boundary point of $\partial D$ (obviously $D$ has a Stein neighborhood basis). This gives a covering of $\partial D$ with a finite number of balls $B_{j}$, mappings $\Phi_{j} \in \mathcal{O}\left(\bar{D}, \mathbb{C}^{n}\right)$ and strongly convex $\mathcal{C}^{\infty}$-smooth domains $C_{j}, j=1, \ldots, N$, such that

- $\Phi_{j}(D) \subset C_{j}$;
- $\Phi_{j}(\bar{D}) \subset \bar{C}_{j}$;
- $\Phi_{j}\left(B_{j} \backslash \bar{D}\right) \subset \mathbb{C}^{n} \backslash \bar{C}_{j}$;
- $\Phi_{j}^{-1}\left(\Phi_{j}\left(B_{j}\right)\right)=B_{j}$;
- $\left.\Phi_{j}\right|_{B_{j}}: B_{j} \rightarrow \Phi_{j}\left(B_{j}\right)$ is biholomorphic.

Therefore, one may choose $c>0$ such that every $C_{j}$ satisfies the exterior ball condition with $c$, i.e. for any $x \in \partial C_{j}$ there is a ball of radius $c$ containing $C_{j}$ and tangent to $\partial C_{j}$ at $x$, every ball of radius $1 / c$ intersecting $\partial D$ is contained in some $B_{j}$ (here one may use a standard argument invoking the Lebesgue number) and the conditions (c), (d) are also satisfied (with $\Phi:=\Phi_{j}$ ).

In this section we use the words 'uniform', 'uniformly' if $(D, z) \in \mathcal{D}(c)$. This means that estimates will depend only on $c$ and will be independent of $D$ and $z$ if $(D, z) \in \mathcal{D}(c)$, and of $E$-mappings of $D$ mapping 0 to $z$ (which we denote by $f:(\mathbb{D}, 0) \rightarrow(D, z))$.

Now we go back to the situation when $D \subset \mathbb{C}^{n}, n \geq 2$, is a bounded strongly linearly convex domain with real-analytic boundary.

Proposition 4.5. Let $f:(\mathbb{D}, 0) \rightarrow(D, z)$ be an E-mapping. Then

$$
\operatorname{dist}(f(\zeta), \partial D) \leq C(1-|\zeta|), \quad \zeta \in \overline{\mathbb{D}}
$$

with $C>0$ uniform if $(D, z) \in \mathcal{D}(c)$.
Proof. There exists a uniform $C_{1}$ such that

$$
\text { if } \operatorname{dist}(w, \partial D) \geq 1 / c \text { then } k_{D}(w, z)<C_{1}
$$

Indeed, let $\operatorname{dist}(w, \partial D) \geq 1 / c$ and let balls $B_{0}, \ldots, B_{m}$ with centers $b_{0}, \ldots, b_{m}$
be chosen for the points $w, z$ as in the condition (3) of Definition 4.1. Then

$$
\begin{aligned}
k_{D}(w, z) & \leq k_{D}\left(w, b_{0}\right)+\sum_{j=0}^{m-1} k_{D}\left(b_{j}, b_{j+1}\right)+k_{D}\left(b_{m}, z\right) \\
& \leq k_{B_{n}(w, 1 / c)}\left(w, b_{0}\right)+\sum_{j=0}^{m-1} k_{B_{j}}\left(b_{j}, b_{j+1}\right)+k_{B_{n}(z, 1 / c)}\left(b_{m}, z\right) \\
& =p\left(0, \frac{\left|w-b_{0}\right|}{1 / c}\right)+\sum_{j=0}^{m-1} p\left(0, \frac{\left|b_{j}-b_{j+1}\right|}{1 /(2 c)}\right)+p\left(0, \frac{\left|b_{m}-z\right|}{1 / c}\right) \\
& \leq(m+2) p(0,1 / 2) \leq\left(8 c^{2}+2\right) p(0,1 / 2)=: C_{1}
\end{aligned}
$$

If $\zeta \in \mathbb{D}$ is such that $\operatorname{dist}(f(\zeta), \partial D) \geq 1 / c$ then

$$
k_{D}(f(0), f(\zeta)) \leq C_{2}-\frac{1}{2} \log \operatorname{dist}(f(\zeta), \partial D)
$$

with a uniform $C_{2}:=C_{1}+\frac{1}{2} \log c$.
In the other case, i.e. if $\operatorname{dist}(f(\zeta), \partial D)<1 / c$, denote by $\eta$ the nearest point to $f(\zeta)$ lying on $\partial D$. Let $w \in D$ be the center of a ball $B$ of radius $1 / c$ tangent to $\partial D$ at $\eta$. By the condition (2) of Definition 4.1 we have $B \subset D$. Hence

$$
\begin{aligned}
& k_{D}(f(0), f(\zeta)) \leq k_{D}(f(0), w)+k_{D}(w, f(\zeta)) \\
& \leq C_{1}+k_{B}(w, f(\zeta)) \leq C_{1}+\frac{1}{2} \log 2-\frac{1}{2} \log \left(1-\frac{|f(\zeta)-w|}{1 / c}\right) \\
& =C_{1}+\frac{1}{2} \log 2-\frac{1}{2} \log (c \operatorname{dist}(f(\zeta), \partial B))=C_{3}-\frac{1}{2} \log \operatorname{dist}(f(\zeta), \partial D)
\end{aligned}
$$

with a uniform $C_{3}:=C_{1}+\frac{1}{2} \log (2 / c)$.
We have obtained the same type estimates in both cases. On the other hand, by Corollary 3.1.2,

$$
k_{D}(f(0), f(\zeta))=p(0, \zeta) \geq-\frac{1}{2} \log (1-|\zeta|)
$$

which finishes the proof.
Recall that we have assumed that $\rho$ is of the form (7).
Proposition 4.6. Let $f:(\mathbb{D}, 0) \rightarrow(D, z)$ be an E-mapping. Then

$$
C_{1}<\rho(\zeta)^{-1}<C_{2}, \quad \zeta \in \mathbb{T}
$$

where $C_{1}, C_{2}$ are uniform if $(D, z) \in \mathcal{D}(c)$.
Proof. For the upper estimate fix $\zeta_{0} \in \mathbb{T}$. Set $B:=B_{n}\left(f\left(\zeta_{0}\right), 1 / c\right)$ and let $\Phi \in \mathcal{O}\left(\bar{D}, \mathbb{C}^{n}\right)$ be as in the condition (4) of Definition 4.1 for $B$. One can assume that $f\left(\zeta_{0}\right)=\Phi\left(f\left(\zeta_{0}\right)\right)=0$ and $\nu_{D}(0)=\nu_{\Phi(D)}(0)=(1,0, \ldots, 0)$.

Then $\Phi(D)$ is contained in the half-space $\left\{w \in \mathbb{C}^{n}: \operatorname{Re} w_{1}<0\right\}$. Putting $h:=\Phi \circ f$ we have

$$
h_{1}(\mathbb{D}) \subset\left\{w_{1} \in \mathbb{C}: \operatorname{Re} w_{1}<0\right\} .
$$

By the Schwarz Lemma on the half-plane,

$$
\begin{equation*}
\left|h_{1}^{\prime}\left(t \zeta_{0}\right)\right| \leq \frac{-2 \operatorname{Re} h_{1}\left(t \zeta_{0}\right)}{1-\left|t \zeta_{0}\right|^{2}} \tag{8}
\end{equation*}
$$

Let $\delta$ be the signed boundary distance of $\Phi(D)$, i.e.

$$
\delta(x):= \begin{cases}-\operatorname{dist}(x, \partial \Phi(D)), & x \in \Phi(D) \\ \operatorname{dist}(x, \partial \Phi(D)), & x \notin \Phi(D)\end{cases}
$$

It is a defining function of $\Phi(D)$ in a neighborhood of 0 (recall that $\Phi^{-1}(\Phi(B))$ $=B$ ). Observe that

$$
\delta(x)=\delta(0)+\operatorname{Re}\langle\nabla \delta(0), x\rangle+O\left(|x|^{2}\right)=\operatorname{Re} x_{1}+O\left(|x|^{2}\right)
$$

If $x \in \Phi(D)$ tends transversally to 0 , then the angle between the vector $x$ and the hyperplane $\left\{w \in \mathbb{C}^{n}: \operatorname{Re} w_{1}=0\right\}$ is separated from 0 , i.e. its sine is $\left(-\operatorname{Re} x_{1}\right) /|x|>\varepsilon$ for some $\varepsilon>0$ independent of $x$. Thus

$$
\frac{\delta(x)}{\operatorname{Re} x_{1}}=1+O(|x|) \quad \text { as } x \rightarrow 0 \text { transversally }
$$

Consequently,

$$
\begin{equation*}
-\operatorname{Re} x_{1} \leq 2 \operatorname{dist}(x, \partial \Phi(D)) \quad \text { as } x \rightarrow 0 \text { transversally. } \tag{9}
\end{equation*}
$$

We know that $t \mapsto f\left(t \zeta_{0}\right)$ hits $\partial D$ transversally. Therefore, $t \mapsto h\left(t \zeta_{0}\right)$ hits $\partial \Phi(D)$ transversally as well. Indeed, we have

$$
\begin{align*}
\left\langle\left.\frac{d}{d t} h\left(t \zeta_{0}\right)\right|_{t=1}\right. & \left., \nu_{\Phi(D)}\left(h\left(\zeta_{0}\right)\right)\right\rangle  \tag{10}\\
& =\left\langle\Phi^{\prime}(0) f^{\prime}\left(\zeta_{0}\right) \zeta_{0}, \frac{\left(\Phi^{-1}\right)^{\prime}(0)^{*} \nabla r(0)}{\left|\left(\Phi^{-1}\right)^{\prime}(0)^{*} \nabla r(0)\right|}\right\rangle \\
& =\frac{\left\langle\zeta_{0} f^{\prime}\left(\zeta_{0}\right), \nabla r(0)\right\rangle}{\mid\left(\Phi^{\prime}(0)^{-1}\right) * \overline{\nabla r(0) \mid}}=\frac{\left\langle\zeta_{0} f^{\prime}\left(\zeta_{0}\right), \nu_{D}\left(f\left(\zeta_{0}\right)\right)\right| \nabla r(0)| \rangle}{\left|\left(\Phi^{\prime}(0)^{-1}\right) * \nabla r(0)\right|}
\end{align*}
$$

where $r$ is a defining function of $D$. In particular,

$$
\operatorname{Re}\left\langle\left.\frac{d}{d t} h\left(t \zeta_{0}\right)\right|_{t=1}, \nu_{\Phi(D)}\left(h\left(\zeta_{0}\right)\right)\right\rangle=\frac{\rho\left(\zeta_{0}\right)^{-1}|\nabla r(0)|}{\left|\left(\Phi^{\prime}(0)^{-1}\right) * \overline{\nabla r(0)}\right|} \neq 0
$$

This proves that $t \mapsto h\left(t \zeta_{0}\right)$ hits $\partial \Phi(D)$ transversally.
Consequently, we may put $x=h\left(t \zeta_{0}\right)$ into (9) to get

$$
\begin{equation*}
\frac{-2 \operatorname{Re} h_{1}\left(t \zeta_{0}\right)}{1-\left|t \zeta_{0}\right|^{2}} \leq \frac{4 \operatorname{dist}\left(h\left(t \zeta_{0}\right), \partial \Phi(D)\right)}{1-\left|t \zeta_{0}\right|^{2}}, \quad t \rightarrow 1 \tag{11}
\end{equation*}
$$

But $\Phi$ is a biholomorphism near 0 , so

$$
\begin{equation*}
\frac{4 \operatorname{dist}\left(h\left(t \zeta_{0}\right), \partial \Phi(D)\right)}{1-\left|t \zeta_{0}\right|^{2}} \leq C_{3} \frac{\operatorname{dist}\left(f\left(t \zeta_{0}\right), \partial D\right)}{1-\left|t \zeta_{0}\right|}, \quad t \rightarrow 1, \tag{12}
\end{equation*}
$$

where $C_{3}$ is uniform (thanks to the condition (4)(c) of Definition 4.1). By Proposition 4.5, the term on the right side of (12) does not exceed some uniform constant.

It follows from 10) that

$$
\begin{aligned}
\rho\left(\zeta_{0}\right)^{-1} & =\left|\left\langle f^{\prime}\left(\zeta_{0}\right) \zeta_{0}, \nu_{D}\left(f\left(\zeta_{0}\right)\right)\right\rangle\right| \\
& \leq C_{4}\left|\left\langle h^{\prime}\left(\zeta_{0}\right), \nu_{\Phi(D)}\left(h\left(\zeta_{0}\right)\right)\right\rangle\right|=C_{4}\left|h_{1}^{\prime}\left(\zeta_{0}\right)\right|=\lim _{t \rightarrow 1} C_{4}\left|h_{1}^{\prime}\left(t \zeta_{0}\right)\right|
\end{aligned}
$$

with a uniform $C_{4}$ (here we use the condition (4)(c) of Definition 4.1 again). Combining (8), (11) and (12) we get the upper estimate for $\rho\left(\zeta_{0}\right)^{-1}$.

Now we prove the lower estimate. Let $r$ be the signed boundary distance to $\partial D$. For $\varepsilon=1 / c$ the function

$$
\varrho(w):=-\log (\varepsilon-r(w))+\log \varepsilon, \quad w \in D_{\varepsilon},
$$

where $D_{\varepsilon}$ is the $\varepsilon$-hull of $D$, is plurisubharmonic and defining for $D$. Indeed, we have

$$
-\log (\varepsilon-r(w))=-\log \operatorname{dist}\left(w, \partial D_{\varepsilon}\right), \quad w \in D_{\varepsilon},
$$

and $D_{\varepsilon}$ is pseudoconvex.
Therefore, the function

$$
v:=\varrho \circ f: \overline{\mathbb{D}} \rightarrow(-\infty, 0]
$$

is subharmonic on $\mathbb{D}$. Moreover, since $f$ maps $\mathbb{T}$ in $\partial D$, we infer that $v=0$ on $\mathbb{T}$. Further, since $|f(\lambda)-z|<c$ for $\lambda \in \mathbb{D}$, we have

$$
|f(\lambda)-z|<\frac{1}{2 c} \quad \text { if }|\lambda| \leq \frac{1}{2 c^{2}} .
$$

Therefore, for a fixed $\zeta_{0} \in \mathbb{T}$,
$M_{\zeta_{0}}(x):=\max _{t \in[0,2 \pi]} v\left(\zeta_{0} e^{x+i t}\right) \leq-\log \left(1+\frac{1}{2 c \varepsilon}\right)=:-C_{5} \quad$ if $x \leq-\log \left(2 c^{2}\right)$.
Since $M_{\zeta_{0}}$ is convex for $x \leq 0$ and $M_{\zeta_{0}}(0)=0$, we get

$$
v\left(\zeta_{0} e^{x}\right) \leq M_{\zeta_{0}}(x) \leq \frac{C_{5} x}{\log \left(2 c^{2}\right)} \quad \text { for }-\log \left(2 c^{2}\right) \leq x \leq 0 .
$$

Hence (remember that $v\left(\zeta_{0}\right)=0$ )

$$
\begin{align*}
\frac{C_{5}}{\log \left(2 c^{2}\right)} & \leq\left.\frac{d}{d x} v\left(\zeta_{0} e^{x}\right)\right|_{x=0}=\sum_{j=1}^{n} \frac{\partial \varrho}{\partial z_{j}}\left(f\left(\zeta_{0}\right)\right) f_{j}^{\prime}\left(\zeta_{0}\right) \zeta_{0}  \tag{13}\\
& =\left\langle\zeta_{0} f^{\prime}\left(\zeta_{0}\right), \nabla \varrho\left(f\left(\zeta_{0}\right)\right)\right\rangle=\rho\left(\zeta_{0}\right)^{-1}\left|\nabla \varrho\left(f\left(\zeta_{0}\right)\right)\right|
\end{align*}
$$

Moreover,

$$
\begin{aligned}
\left|\nabla \varrho\left(f\left(\zeta_{0}\right)\right)\right| & =\left\langle\nabla \varrho\left(f\left(\zeta_{0}\right)\right), \frac{\nabla \varrho\left(f\left(\zeta_{0}\right)\right)}{\left|\nabla \varrho\left(f\left(\zeta_{0}\right)\right)\right|}\right\rangle_{\mathbb{R}}=\left\langle\nabla \varrho\left(f\left(\zeta_{0}\right)\right), \nu_{D}\left(f\left(\zeta_{0}\right)\right)\right\rangle_{\mathbb{R}} \\
& =\frac{\partial \varrho}{\partial \nu_{D}}\left(f\left(\zeta_{0}\right)\right)=\lim _{t \rightarrow 0} \frac{\varrho\left(f\left(\zeta_{0}\right)+t \nu_{D}\left(f\left(\zeta_{0}\right)\right)\right)-\varrho\left(f\left(\zeta_{0}\right)\right)}{t}=\frac{1}{\varepsilon}=c,
\end{aligned}
$$

as $r(a+t \nu(a))=t$ if $a \in \partial D$ and $t \in \mathbb{R}$ is small enough. This, together with (13), finishes the proof of the lower estimate.

Proposition 4.7. Let $f:(\mathbb{D}, 0) \rightarrow(D, z)$ be an E-mapping. Then

$$
\left|f\left(\zeta_{1}\right)-f\left(\zeta_{2}\right)\right| \leq C \sqrt{\left|\zeta_{1}-\zeta_{2}\right|}, \quad \zeta_{1}, \zeta_{2} \in \overline{\mathbb{D}}
$$

where $C$ is uniform if $(D, z) \in \mathcal{D}(c)$.
Proof. Let $\zeta_{0} \in \mathbb{D}$ be such that $1-\left|\zeta_{0}\right|<1 /(c C)$, where $C$ is as in Proposition 4.5. Then $B:=B_{n}\left(f\left(\zeta_{0}\right), 1 / c\right)$ intersects $\partial D$. Take $\Phi$ for the ball $B$ from the condition (4) of Definition 4.1. Let $w$ denote the nearest point to $\Phi\left(f\left(\zeta_{0}\right)\right)$ lying on $\partial \Phi(D)$. From the conditions (4)(b)-(c) of Definition 4.1 we find that there is a uniform constant $r<1$ such that the point $w$ belongs to $\Phi(B \cap \partial D)$ provided that $\left|\zeta_{0}\right| \geq r$.

From the condition (4) (a) of Definition 4.1 we deduce that there is $w_{0}$ such that $\Phi(D) \subset B_{n}\left(w_{0}, c\right)$ and the ball $B_{n}\left(w_{0}, c\right)$ is tangent to $\partial \Phi(D)$ at $w$. Let

$$
h(\zeta):=(\Phi \circ f)\left(\frac{\zeta_{0}-\zeta}{1-\overline{\zeta_{0}} \zeta}\right), \quad \zeta \in \mathbb{D}
$$

Then $h$ is holomorphic, $h(\mathbb{D}) \subset B_{n}\left(w_{0}, c\right)$ and $h(0)=\Phi\left(f\left(\zeta_{0}\right)\right)$. Using the Schwarz Lemma for the unit ball, we get

$$
\begin{aligned}
\left|h^{\prime}(0)\right| & \leq \sqrt{c^{2}-\left|h(0)-w_{0}\right|^{2}} \\
& \leq \sqrt{2 c\left(c-\left|\Phi\left(f\left(\zeta_{0}\right)\right)-w_{0}\right|\right)}=\sqrt{2 c\left(\left|w_{0}-w\right|-\left|\Phi\left(f\left(\zeta_{0}\right)\right)-w_{0}\right|\right)} \\
& \leq \sqrt{2 c} \sqrt{\left|\Phi\left(f\left(\zeta_{0}\right)\right)-w\right|}=\sqrt{2 c} \sqrt{\operatorname{dist}\left(\Phi\left(f\left(\zeta_{0}\right)\right), \partial \Phi(D)\right)} .
\end{aligned}
$$

Since

$$
h^{\prime}(0)=\left.\Phi^{\prime}\left(f\left(\zeta_{0}\right)\right) f^{\prime}\left(\zeta_{0}\right) \frac{d}{d \zeta} \frac{\zeta_{0}-\zeta}{1-\overline{\zeta_{0}} \zeta}\right|_{\zeta=0},
$$

we deduce by the condition $(4)(c)$ of Definition 4.1 that

$$
\left|h^{\prime}(0)\right| \geq C_{1}\left|f^{\prime}\left(\zeta_{0}\right)\right|\left(1-\left|\zeta_{0}\right|^{2}\right)
$$

with a uniform $C_{1}$. Further,

$$
\begin{aligned}
\left|f^{\prime}\left(\zeta_{0}\right)\right| & \leq \frac{\left|h^{\prime}(0)\right|}{C_{1}\left(1-\left|\zeta_{0}\right|^{2}\right)} \leq \frac{\sqrt{2 c}}{C_{1}} \frac{\sqrt{\operatorname{dist}\left(\Phi\left(f\left(\zeta_{0}\right)\right), \partial \Phi(D)\right)}}{1-\left|\zeta_{0}\right|^{2}} \\
& \leq C_{2} \frac{\sqrt{\operatorname{dist}\left(f\left(\zeta_{0}\right), \partial D\right)}}{1-\left|\zeta_{0}\right|^{2}}
\end{aligned}
$$

where $C_{2}$ is uniform. Combining this with Proposition 4.5 we find that

$$
\begin{equation*}
\left|f^{\prime}\left(\zeta_{0}\right)\right| \leq C_{3} \frac{\sqrt{1-\left|\zeta_{0}\right|}}{1-\left|\zeta_{0}\right|^{2}}=\frac{C_{3}}{\sqrt{1-\left|\zeta_{0}\right|}} \tag{14}
\end{equation*}
$$

where $C_{3}$ is a uniform constant.
We have shown that (14) holds for $r \leq\left|\zeta_{0}\right|<1$ with a uniform $r<1$. For $\left|\zeta_{0}\right|<r$ we estimate in the following way:

$$
\left|f^{\prime}\left(\zeta_{0}\right)\right| \leq \max _{|\zeta|=r}\left|f^{\prime}(\zeta)\right| \leq \frac{C_{3}}{\sqrt{1-r}} \leq \frac{C_{4}}{\sqrt{1-\left|\zeta_{0}\right|}}
$$

with a uniform $C_{4}:=C_{3} / \sqrt{1-r}$.
Using Theorems 8.2 .1 and 8.2 .2 with $\alpha=1 / 2$, we finish the proof.
Proposition 4.8. Let $f:(\mathbb{D}, 0) \rightarrow(D, z)$ be an E-mapping. Then

$$
\left|\rho\left(\zeta_{1}\right)-\rho\left(\zeta_{2}\right)\right| \leq C \sqrt{\left|\zeta_{1}-\zeta_{2}\right|}, \quad \zeta_{1}, \zeta_{2} \in \mathbb{T}
$$

where $C$ is uniform if $(D, z) \in \mathcal{D}(c)$.
Proof. It suffices to prove that there exist uniform $C, C_{1}>0$ such that

$$
\left|\rho\left(\zeta_{1}\right)-\rho\left(\zeta_{2}\right)\right| \leq C \sqrt{\left|\zeta_{1}-\zeta_{2}\right|}, \quad \zeta_{1}, \zeta_{2} \in \mathbb{T},\left|\zeta_{1}-\zeta_{2}\right|<C_{1}
$$

Fix $\zeta_{1} \in \mathbb{T}$. Without loss of generality we may assume that $\nu_{D, 1}\left(f\left(\zeta_{1}\right)\right)$ $=1$. Let $0<C_{1} \leq 1 / 4$ be uniform and such that

$$
\left|\nu_{D, 1}(f(\zeta))-1\right|<1 / 2, \quad \zeta \in \mathbb{T} \cap B_{n}\left(\zeta_{1}, 3 C_{1}\right)
$$

This is possible, since by Proposition 4.7,

$$
\left|\nu_{D}(f(\zeta))-\nu_{D}\left(f\left(\zeta^{\prime}\right)\right)\right| \leq 2 c\left|f(\zeta)-f\left(\zeta^{\prime}\right)\right| \leq C^{\prime} \sqrt{\left|\zeta-\zeta^{\prime}\right|}, \quad \zeta, \zeta^{\prime} \in \mathbb{T}
$$

with a uniform $C^{\prime}>0$. There exists a function $\psi \in \mathcal{C}^{1}(\mathbb{T},[0,1])$ such that $\psi=1$ on $\mathbb{T} \cap B_{n}\left(\zeta_{1}, 2 C_{1}\right)$ and $\psi=0$ on $\mathbb{T} \backslash B_{n}\left(\zeta_{1}, 3 C_{1}\right)$. Then the function $\varphi: \mathbb{T} \rightarrow \mathbb{C}$ defined by

$$
\varphi:=\left(\overline{\nu_{D, 1} \circ f}-1\right) \psi+1
$$

satisfies

- $\varphi(\zeta)=\overline{\nu_{D, 1}(f(\zeta))}, \zeta \in \mathbb{T} \cap B_{n}\left(\zeta_{1}, 2 C_{1}\right) ;$
- $|\varphi(\zeta)-1|<1 / 2, \zeta \in \mathbb{T}$;
- $\varphi$ is uniformly $1 / 2$-Hölder continuous on $\mathbb{T}$, i.e. it is $1 / 2$-Hölder continuous with a uniform constant (remember that $\psi$ was chosen uniformly).
First observe that $\log \varphi$ is well-defined. Using properties listed above, we deduce that $\log \varphi$ and $\operatorname{Im} \log \varphi$ are uniformly $1 / 2$-Hölder continuous on $\mathbb{T}$ as well. The function $\operatorname{Im} \log \varphi$ extends continuously to a function $v: \overline{\mathbb{D}} \rightarrow \mathbb{R}$, harmonic in $\mathbb{D}$. There is a function $h \in \mathcal{O}(\mathbb{D})$ such that $v=\operatorname{Im} h$ in $\mathbb{D}$. Taking $h-\operatorname{Re} h(0)$ instead of $h$, one can assume that $\operatorname{Re} h(0)=0$. By Theorem8.2.3 applied to $i h$, we see that the function $h$ extends continuously on $\overline{\mathbb{D}}$ and $h$ is uniformly $1 / 2$-Hölder continuous on $\overline{\mathbb{D}}$. Hence the function
$u:=\operatorname{Re} h: \overline{\mathbb{D}} \rightarrow \mathbb{R}$ is uniformly $1 / 2$-Hölder continuous on $\overline{\mathbb{D}}$ with a uniform constant $C_{2}$. Furthermore, $u$ is uniformly bounded on $\overline{\mathbb{D}}$, since

$$
|u(\zeta)|=|u(\zeta)-u(0)| \leq C_{2} \sqrt{|\zeta|}, \quad \zeta \in \overline{\mathbb{D}} .
$$

Let $g(\zeta):=\widetilde{f}_{1}(\zeta) e^{-h(\zeta)}$ and $G(\zeta):=g(\zeta) / \zeta$. Then $g \in \mathcal{O}(\mathbb{D}) \cap \mathcal{C}(\overline{\mathbb{D}})$ and $G \in \mathcal{O}\left(\mathbb{D}_{*}\right) \cap \mathcal{C}\left((\overline{\mathbb{D}})_{*}\right)$. Note that, for $\zeta \in \mathbb{T}$,

$$
|g(\zeta)|=\left|\zeta \rho(\zeta) \overline{\nu_{D, 1}(f(\zeta))} e^{-h(\zeta)}\right| \leq \rho(\zeta) e^{-u(\zeta)}
$$

which, combined with Proposition 4.6, the uniform boundedness of $u$ and the maximum principle, gives uniform boundedness of $g$ in $\overline{\mathbb{D}}$. The function $G$ is uniformly bounded in $\overline{\mathbb{D}} \cap B_{n}\left(\zeta_{1}, 2 C_{1}\right)$. Moreover, for $\zeta \in \mathbb{T} \cap B_{n}\left(\zeta_{1}, 2 C_{1}\right)$,

$$
\begin{aligned}
G(\zeta) & =\rho(\zeta) \overline{\nu_{D, 1}(f(\zeta))} e^{-u(\zeta)-i \operatorname{Im} \log \varphi(\zeta)} \\
& =\rho(\zeta) \overline{\nu_{D, 1}(f(\zeta))} e^{-u(\zeta)+\operatorname{Re} \log \varphi(\zeta)} e^{-\log \varphi(\zeta)}=\rho(\zeta) e^{-u(\zeta)+\operatorname{Re} \log \varphi(\zeta)} \in \mathbb{R}
\end{aligned}
$$

By the Reflection Principle one can extend $G$ holomorphically past $\mathbb{T} \cap$ $B_{n}\left(\zeta_{1}, 2 C_{1}\right)$ to a function (denoted by the same letter) uniformly bounded in $B_{n}\left(\zeta_{1}, 2 C_{2}\right)$, where the constant $C_{2}$ is uniform. Hence, from the Cauchy formula, $G$ is uniformly Lipschitz continuous in $B_{n}\left(\zeta_{1}, C_{2}\right)$, and consequently uniformly $1 / 2$-Hölder continuous in $B_{n}\left(\zeta_{1}, C_{2}\right)$.

Finally, the functions $G, h, \nu_{D, 1} \circ f$ are uniformly $1 / 2$-Hölder continuous on $\mathbb{T} \cap B_{n}\left(\zeta_{1}, C_{2}\right)$ and $\left|\nu_{D, 1} \circ f\right|>1 / 2$ on $\mathbb{T} \cap B_{n}\left(\zeta_{1}, C_{2}\right)$, so the function $\rho=G e^{h} / \overline{\nu_{D, 1} \circ f}$ is uniformly $1 / 2$-Hölder continuous on $\mathbb{T} \cap B_{n}\left(\zeta_{1}, C_{2}\right)$.

Proposition 4.9. Let $f:(\mathbb{D}, 0) \rightarrow(D, z)$ be an E-mapping. Then

$$
\left|\tilde{f}\left(\zeta_{1}\right)-\tilde{f}\left(\zeta_{2}\right)\right| \leq C \sqrt{\left|\zeta_{1}-\zeta_{2}\right|}, \quad \zeta_{1}, \zeta_{2} \in \overline{\mathbb{D}}
$$

where $C$ is uniform if $(D, z) \in \mathcal{D}(c)$.
Proof. By Propositions 4.7 and 4.8 we obtain the desired inequality for $\zeta_{1}, \zeta_{2} \in \mathbb{T}$. Theorem 8.2.2 finishes the proof.
5. Openness of the set of $E$-mappings. We shall show that perturbing a little a domain $D$ equipped with an $E$-mapping, we obtain a domain which also has an $E$-mapping close to a given one.

### 5.1. Preliminary results

Proposition 5.1.1. Let $f: \mathbb{D} \rightarrow D$ be an E-mapping. Then there exist domains $G, \widetilde{D}, \widetilde{G} \subset \mathbb{C}^{n}$ and a biholomorphism $\Phi: \widetilde{D} \rightarrow \widetilde{G}$ such that
(1) $\widetilde{D}, \widetilde{G}$ are neighborhoods of $\bar{D}, \bar{G}$ respectively;
(2) $\Phi(D)=G$;
(3) $g(\zeta):=\Phi(f(\zeta))=(\zeta, 0, \ldots, 0), \zeta \in \overline{\mathbb{D}}$;
(4) $\nu_{G}(g(\zeta))=(\zeta, 0, \ldots, 0), \zeta \in \mathbb{T}$;
(5) for any $\zeta \in \mathbb{T}$, a point $g(\zeta)$ is a point of strong linear convexity of $G$.

Proof. Let $U, V$ be the sets from Proposition 3.2.4. We claim that after a linear change of coordinates one can assume that $f_{1}, f_{2}$ do not have common zeroes in $V$.

Since $f^{\prime} \bullet \tilde{f}=1$, at least one of the functions $\tilde{f}_{1}, \ldots, \tilde{f}_{n}$, say $\widetilde{f}_{1}$, is not identically equal to 0 . Let $\lambda_{1}, \ldots, \lambda_{m}$ be all zeroes of $\tilde{f}_{1}$ in $V$. We may find $\alpha \in \mathbb{C}^{n}$ such that

$$
\left(\alpha_{1} \widetilde{f}_{1}+\cdots+\alpha_{n} \widetilde{f}_{n}\right)\left(\lambda_{j}\right) \neq 0, \quad j=1, \ldots, m
$$

Otherwise, for any $\alpha \in \mathbb{C}^{n}$ there would exist $j \in\{1, \ldots, m\}$ such that $\alpha \bullet \widetilde{f}\left(\lambda_{j}\right)=0$, hence

$$
\mathbb{C}^{n}=\bigcup_{j=1}^{m}\left\{\alpha \in \mathbb{C}^{n}: \alpha \bullet \tilde{f}\left(\lambda_{j}\right)=0\right\}
$$

The sets $\left\{\alpha \in \mathbb{C}^{n}: \alpha \bullet \widetilde{f}\left(\lambda_{j}\right)=0\right\}, j=1, \ldots, m$, are $(n-1)$-dimensional complex hyperplanes, so their finite union cannot be the space $\mathbb{C}^{n}$.

Of course, at least one of the numbers $\alpha_{2}, \ldots, \alpha_{n}$, say $\alpha_{2}$, is non-zero. Let

$$
A:=\left[\begin{array}{ccccc}
1 & 0 & 0 & \cdots & 0 \\
\alpha_{1} & \alpha_{2} & \alpha_{3} & \cdots & \alpha_{n} \\
0 & 0 & 1 & \cdots & 0 \\
\vdots & \vdots & \vdots & \ddots & \vdots \\
0 & 0 & 0 & \cdots & 1
\end{array}\right], \quad B:=\left(A^{T}\right)^{-1}
$$

We claim that $B$ is a change of coordinates we are looking for. If $r$ is a defining function of $D$ then $r \circ B^{-1}$ is a defining function of $B(D)$, so $B(D)$ is a bounded strongly linearly convex domain with real-analytic boundary. Let us check that $B f$ is an $E$-mapping of $B(D)$ with associated mappings

$$
\begin{equation*}
A \widetilde{f} \in \mathcal{O}(\overline{\mathbb{D}}) \quad \text { and } \quad \rho \frac{|A \overline{\nabla r \circ f}|}{|\nabla r \circ f|} \in \mathcal{C}^{\omega}(\mathbb{T}) \tag{15}
\end{equation*}
$$

The conditions (1) and (2) of Definition 1.11 are clear. For $\zeta \in \mathbb{T}$ we have

$$
\begin{align*}
\overline{\nu_{B(D)}(B f(\zeta))} & =\frac{\overline{\nabla\left(r \circ B^{-1}\right)(B f(\zeta))}}{\left|\nabla\left(r \circ B^{-1}\right)(B f(\zeta))\right|}  \tag{16}\\
& =\frac{\left(B^{-1}\right)^{T} \overline{\nabla r(f(\zeta))}}{\left|\left(B^{-1}\right)^{T} \overline{\nabla r(f(\zeta))}\right|}=\frac{A \overline{\nabla r(f(\zeta))}}{\mid A \overline{\nabla r(f(\zeta)) \mid}},
\end{align*}
$$

so

Moreover, for $\zeta \in \mathbb{T}, z \in D$,

$$
\begin{aligned}
& \left\langle B z-B f(\zeta), \nu_{B(D)}(B f(\zeta))\right\rangle={\overline{\nu_{B(D)}(B f(\zeta))}}^{T}(B z-B f(\zeta)) \\
& =\frac{\overline{\nabla r(f(\zeta))}^{T} B^{-1} B(z-f(\zeta))}{\left|\left(B^{-1}\right)^{T} \overline{\nabla r(f(\zeta))}\right|}=\frac{|\nabla r(f(\zeta))|}{\left|\left(B^{-1}\right)^{T} \overline{\nabla r(f(\zeta))}\right|}{\overline{\nu_{D}(f(\zeta))}}^{T}(z-f(\zeta)) \\
& =\frac{|\nabla r(f(\zeta))|}{\left|\left(B^{-1}\right)^{T} \overline{\nabla r(f(\zeta))}\right|}\left\langle z-f(\zeta), \nu_{D}(f(\zeta))\right\rangle .
\end{aligned}
$$

Therefore, $B$ is the desired linear change of coordinates, as claimed.
If necessary, we shrink the sets $U, V$ associated with $f$ to sets associated with $B f$. There exist holomorphic functions $h_{1}, h_{2}: V \rightarrow \mathbb{C}$ such that

$$
h_{1} \widetilde{f}_{1}+h_{2} \widetilde{f}_{2} \equiv 1 \quad \text { in } V
$$

Generally, this is a well-known fact for functions on pseudoconvex domains, however in the present case it may be shown elementarily. Indeed, if $\widetilde{f}_{1} \equiv 0$ or $\widetilde{f}_{2} \equiv 0$ then it is obvious. In the opposite case, let $\widetilde{f}_{j}=F_{j} P_{j}, j=1,2$, where $F_{j}$ are holomorphic, non-zero in $V$, and $P_{j}$ are polynomials with all (finitely many) zeroes in $V$. Then $P_{j}$ are relatively prime, so there are polynomials $Q_{j}, j=1,2$, such that

$$
Q_{1} P_{1}+Q_{2} P_{2} \equiv 1
$$

Hence

$$
\frac{Q_{1}}{F_{1}} \widetilde{f}_{1}+\frac{Q_{2}}{F_{2}} \widetilde{f}_{2} \equiv 1 \quad \text { in } V
$$

Consider the mapping $\Psi: V \times \mathbb{C}^{n-1} \rightarrow \mathbb{C}^{n}$ given by

$$
\begin{align*}
& \Psi_{1}(Z):=f_{1}\left(Z_{1}\right)-Z_{2} \widetilde{f}_{2}\left(Z_{1}\right)-h_{1}\left(Z_{1}\right) \sum_{j=3}^{n} Z_{j} \widetilde{f}_{j}\left(Z_{1}\right),  \tag{18}\\
& \Psi_{2}(Z):=f_{2}\left(Z_{1}\right)+Z_{2} \widetilde{f}_{1}\left(Z_{1}\right)-h_{2}\left(Z_{1}\right) \sum_{j=3}^{n} Z_{j} \widetilde{f}_{j}\left(Z_{1}\right), \tag{19}
\end{align*}
$$

We claim that $\Psi$ is biholomorphic in $\Psi^{-1}(U)$. First of all, observe that $\Psi^{-1}(\{z\}) \neq \emptyset$ for any $z \in U$. Indeed, by Proposition 3.2.4 there exists (exactly one) $Z_{1} \in V$ such that

$$
\left(z-f\left(Z_{1}\right)\right) \bullet \tilde{f}\left(Z_{1}\right)=0
$$

The numbers $Z_{j} \in \mathbb{C}, j=3, \ldots, n$, are uniquely determined by the equations

$$
Z_{j}=z_{j}-f_{j}\left(Z_{1}\right)
$$

At least one of the numbers $\widetilde{f}_{1}\left(Z_{1}\right), \widetilde{f}_{2}\left(Z_{1}\right)$, say $\widetilde{f}_{1}\left(Z_{1}\right)$, is non-zero. Let

$$
Z_{2}:=\frac{z_{2}-f_{2}\left(Z_{1}\right)+h_{2}\left(Z_{1}\right) \sum_{j=3}^{n} Z_{j} \tilde{f}_{j}\left(Z_{1}\right)}{\widetilde{f}_{1}\left(Z_{1}\right)}
$$

Then we easily check that the equality

$$
z_{1}=f_{1}\left(Z_{1}\right)-Z_{2} \widetilde{f}_{2}\left(Z_{1}\right)-h_{1}\left(Z_{1}\right) \sum_{j=3}^{n} Z_{j} \widetilde{f}_{j}\left(Z_{1}\right)
$$

is equivalent to $\left(z-f\left(Z_{1}\right)\right) \bullet \tilde{f}\left(Z_{1}\right)=0$, which is true.
To finish the proof of biholomorphicity of $\Psi$ in $\Psi^{-1}(U)$ it suffices to check that $\Psi$ is injective in $\Psi^{-1}(U)$. Let us take $Z, W$ such that $\Psi(Z)=\Psi(W)=$ $z \in U$. By a direct computation both $\zeta=Z_{1} \in V$ and $\zeta=W_{1} \in V$ solve the equation

$$
(z-f(\zeta)) \bullet \widetilde{f}(\zeta)=0
$$

From Proposition 3.2.4 we infer that it has exactly one solution. Hence $Z_{1}=$ $W_{1}$. By 20 we have $Z_{j}=W_{j}$ for $j=3, \ldots, n$. Finally $Z_{2}=W_{2}$ follows from one of the equations (18), (19).

Let $G:=\Psi^{-1}(D), \widetilde{D}:=U, \widetilde{G}:=\Psi^{-1}(U), \Phi:=\Psi^{-1}$. We have

$$
\Psi_{j}(\zeta, 0, \ldots, 0)=f_{j}(\zeta), \quad j=1, \ldots, n
$$

so $\Phi(f(\zeta))=(\zeta, 0, \ldots, 0), \zeta \in \overline{\mathbb{D}}$. Put $g(\zeta):=\Phi(f(\zeta)), \zeta \in \overline{\mathbb{D}}$. Note that the entries of the matrix $\Psi^{\prime}(g(\zeta))$ are

$$
\begin{aligned}
\frac{\partial \Psi_{1}}{\partial Z_{1}}(g(\zeta))=f_{1}^{\prime}(\zeta), \frac{\partial \Psi_{1}}{\partial Z_{2}}(g(\zeta))=-\widetilde{f}_{2}(\zeta), & \frac{\partial \Psi_{1}}{\partial Z_{j}}(g(\zeta))=-h_{1}(\zeta) \widetilde{f}_{j}(\zeta), j \geq 3 \\
\frac{\partial \Psi_{2}}{\partial Z_{1}}(g(\zeta))=f_{2}^{\prime}(\zeta), \frac{\partial \Psi_{2}}{\partial Z_{2}}(g(\zeta))=\widetilde{f}_{1}(\zeta), & \frac{\partial \Psi_{2}}{\partial Z_{j}}(g(\zeta))=-h_{2}(\zeta) \widetilde{f}_{j}(\zeta), j \geq 3 \\
\frac{\partial \Psi_{k}}{\partial Z_{1}}(g(\zeta))=f_{k}^{\prime}(\zeta), \frac{\partial \Psi_{k}}{\partial Z_{2}}(g(\zeta))=0, & \frac{\partial \Psi_{k}}{\partial Z_{j}}(g(\zeta))=\delta_{j}^{k}, j, k \geq 3
\end{aligned}
$$

Thus $\Psi^{\prime}(g(\zeta))^{T} \widetilde{f}(\zeta)=(1,0, \ldots, 0), \zeta \in \overline{\mathbb{D}}\left(\right.$ as $\left.f^{\prime} \bullet \widetilde{f}=1\right)$. Let us take a defining function $r$ of $D$. Then $r \circ \Psi$ is a defining function of $G$. Therefore,

$$
\begin{aligned}
& \nu_{G}(g(\zeta))=\frac{\nabla(r \circ \Psi)(g(\zeta))}{|\nabla(r \circ \Psi)(g(\zeta))|}=\frac{{\overline{\Psi^{\prime}(g(\zeta))}}^{T} \nabla r(f(\zeta))}{\left|\overline{\Psi^{\prime}(g(\zeta))^{T}} \nabla r(f(\zeta))\right|} \\
&=\frac{{\overline{\Psi^{\prime}(g(\zeta))}}^{T} \frac{\bar{f}(\zeta)}{\zeta \rho(\zeta)}}{|\nabla r(f(\zeta))|} \\
&\left|{\overline{\Psi^{\prime}(g(\zeta))}}^{T} \frac{\widetilde{f}(\zeta)}{\overline{\zeta \rho(\zeta)}}\right| \nabla r(f(\zeta))|\mid
\end{aligned}=g(\zeta), \quad \zeta \in \mathbb{T} .
$$

It remains to prove the fifth condition. By (5) (see Definition 1.1) we have
to show that

$$
\begin{equation*}
\sum_{j, k=1}^{n} \frac{\partial^{2}(r \circ \Psi)}{\partial z_{j} \partial \bar{z}_{k}}(g(\zeta)) X_{j} \bar{X}_{k}>\left|\sum_{j, k=1}^{n} \frac{\partial^{2}(r \circ \Psi)}{\partial z_{j} \partial z_{k}}(g(\zeta)) X_{j} X_{k}\right| \tag{21}
\end{equation*}
$$

for $\zeta \in \mathbb{T}$ and $X \in\left(\mathbb{C}^{n}\right)_{*}$ with

$$
\sum_{j=1}^{n} \frac{\partial(r \circ \Psi)}{\partial z_{j}}(g(\zeta)) X_{j}=0
$$

i.e. $X_{1}=0$. We have

$$
\begin{aligned}
\sum_{j, k=1}^{n} \frac{\partial^{2}(r \circ \Psi)}{\partial z_{j} \partial \bar{z}_{k}}(g(\zeta)) X_{j} \bar{X}_{k} & =\sum_{j, k, s, t=1}^{n} \frac{\partial^{2} r}{\partial z_{s} \partial \bar{z}_{t}}(f(\zeta)) \frac{\partial \Psi_{s}}{\partial z_{j}}(g(\zeta)) \overline{\frac{\partial \Psi_{t}}{\partial z_{k}}(g(\zeta))} X_{j} \bar{X}_{k} \\
& =\sum_{s, t=1}^{n} \frac{\partial^{2} r}{\partial z_{s} \partial \bar{z}_{t}}(f(\zeta)) Y_{s} \bar{Y}_{t}
\end{aligned}
$$

where

$$
Y:=\Psi^{\prime}(g(\zeta)) X
$$

Note that $Y \neq 0$. Additionally

$$
\sum_{s=1}^{n} \frac{\partial r}{\partial z_{s}}(f(\zeta)) Y_{s}=\sum_{j, s=1}^{n} \frac{\partial r}{\partial z_{s}}(f(\zeta)) \frac{\partial \Psi_{s}}{\partial z_{j}}(g(\zeta)) X_{j}=\sum_{j=1}^{n} \frac{\partial(r \circ \Psi)}{\partial z_{j}}(g(\zeta)) X_{j}=0
$$

Therefore, by the strong linear convexity of $D$ at $f(\zeta)$,

$$
\sum_{s, t=1}^{n} \frac{\partial^{2} r}{\partial z_{s} \partial \bar{z}_{t}}(f(\zeta)) Y_{s} \bar{Y}_{t}>\left|\sum_{s, t=1}^{n} \frac{\partial^{2} r}{\partial z_{s} \partial z_{t}}(f(\zeta)) Y_{s} Y_{t}\right|
$$

To finish the proof observe that

$$
\begin{aligned}
& \left|\sum_{j, k=1}^{n} \frac{\partial^{2}(r \circ \Psi)}{\partial z_{j} \partial z_{k}}(g(\zeta)) X_{j} X_{k}\right| \\
& =\left\lvert\, \sum_{j, k, s, t=1}^{n} \frac{\partial^{2} r}{\partial z_{s} \partial z_{t}}(f(\zeta)) \frac{\partial \Psi_{s}}{\partial z_{j}}(g(\zeta)) \frac{\partial \Psi_{t}}{\partial z_{k}}(g(\zeta)) X_{j} X_{k}\right. \\
& \\
& \left.\quad+\sum_{j, k, s=1}^{n} \frac{\partial r}{\partial z_{s}}(f(\zeta)) \frac{\partial^{2} \Psi_{s}}{\partial z_{j} \partial z_{k}}(g(\zeta)) X_{j} X_{k} \right\rvert\, \\
& =\left|\sum_{s, t=1}^{n} \frac{\partial^{2} r}{\partial z_{s} \partial z_{t}}(f(\zeta)) Y_{s} Y_{t}+\sum_{j, k=2}^{n} \sum_{s=1}^{n} \frac{\partial r}{\partial z_{s}}(f(\zeta)) \frac{\partial^{2} \Psi_{s}}{\partial z_{j} \partial z_{k}}(g(\zeta)) X_{j} X_{k}\right|
\end{aligned}
$$

and

$$
\frac{\partial^{2} \Psi_{s}}{\partial z_{j} \partial z_{k}}(g(\zeta))=0, \quad j, k \geq 2, s \geq 1
$$

which gives (21).

REMARK 5.1.2. Let $D$ be a bounded domain in $\mathbb{C}^{n}$ and let $f: \mathbb{D} \rightarrow D$ be a (weak) stationary mapping such that $\partial D$ is real-analytic in a neighborhood of $f(\mathbb{T})$. Assume moreover that $U$ is a neighborhood of $f(\overline{\mathbb{D}})$, a mapping $\Theta: U \rightarrow \mathbb{C}^{n}$ is biholomorphic onto its image and the set $D \cap U$ is connected. Then $\Theta \circ f$ is a (weak) stationary mapping of $G:=\Theta(D \cap U)$.

In particular, if $U_{1}, U_{2}$ are neighborhoods of the closures of domains $D_{1}, D_{2}$ with real-analytic boundaries and $\Theta: U_{1} \rightarrow U_{2}$ is a biholomorphism such that $\Theta\left(D_{1}\right)=D_{2}$, then $\Theta$ maps (weak) stationary mappings of $D_{1}$ onto (weak) stationary mappings of $D_{2}$.

Proof. Actually, it is clear that the first two conditions of the definition of (weak) stationary mappings are preserved by $\Theta$. To show the third one we proceed as in (15-17). Let $f: \mathbb{D} \rightarrow D$ be a (weak) stationary mapping. The candidates for the mappings in the condition (3) (resp. ( $3^{\prime}$ )) of Definition 1.11 for $\Theta \circ f$ in the domain $G$ are

$$
\left(\left(\Theta^{\prime} \circ f\right)^{-1}\right)^{T} \tilde{f} \quad \text { and } \quad \rho \frac{\left|\left(\left(\Theta^{\prime} \circ f\right)^{-1}\right)^{T} \overline{\nabla r \circ f}\right|}{|\nabla r \circ f|}
$$

Indeed, for $\zeta \in \mathbb{T}$,

$$
\begin{aligned}
\overline{\nu_{G}(\Theta(f(\zeta)))} & =\frac{\overline{\nabla\left(r \circ \Theta^{-1}\right)(\Theta(f(\zeta)))}}{\left|\nabla\left(r \circ \Theta^{-1}\right)(\Theta(f(\zeta)))\right|}=\frac{\left[\left(\Theta^{-1}\right)^{\prime}(\Theta(f(\zeta)))\right]^{T} \overline{\nabla r(f(\zeta))}}{\mid\left[\left(\Theta^{-1}\right)^{\prime}(\Theta(f(\zeta)))\right]^{T} \overline{\nabla r(f(\zeta)) \mid}} \\
& =\frac{\left(\Theta^{\prime}(f(\zeta))^{-1}\right)^{T} \overline{\nabla r(f(\zeta))}}{\left|\left(\Theta^{\prime}(f(\zeta))^{-1}\right)^{T} \overline{\nabla r(f(\zeta))}\right|},
\end{aligned}
$$

hence

$$
\begin{aligned}
\zeta \rho(\zeta) \frac{\left|\left(\Theta^{\prime}(f(\zeta))^{-1}\right)^{T} \overline{\nabla r(f(\zeta))}\right|}{|\nabla r(f(\zeta))|} \overline{\nu_{G}(\Theta(f(\zeta)))} & =\zeta \rho(\zeta)\left(\Theta^{\prime}(f(\zeta))^{-1}\right)^{T} \overline{\nu_{D}(f(\zeta))} \\
& =\left(\Theta^{\prime}(f(\zeta))^{-1}\right)^{T} \widetilde{f}(\zeta)
\end{aligned}
$$

5.2. Situation ( $\dagger$ ). Consider the following situation, denoted by ( $\dagger$ ) (with data $D_{0}$ and $U_{0}$ ):

- $D_{0}$ is a bounded domain in $\mathbb{C}^{n}, n \geq 2$;
- $f_{0}: \overline{\mathbb{D}} \ni \zeta \mapsto(\zeta, 0, \ldots, 0) \in \bar{D}_{0}, \zeta \in \overline{\mathbb{D}} ;$
- $f_{0}(\mathbb{D}) \subset D_{0}$;
- $f_{0}(\mathbb{T}) \subset \partial D_{0}$;
- $\nu_{D_{0}}\left(f_{0}(\zeta)\right)=(\zeta, 0, \ldots, 0), \zeta \in \mathbb{T}$;
- for any $\zeta \in \mathbb{T}$, a point $f_{0}(\zeta)$ is a point of strong linear convexity of $D_{0}$;
- $\partial D_{0}$ is real-analytic in a neighborhood $U_{0}$ of $f_{0}(\mathbb{T})$ with a function $r_{0}$;
- $\left|\nabla r_{0}\right|=1$ on $f_{0}(\mathbb{T})$; in particular, $r_{0 z}\left(f_{0}(\zeta)\right)=(\bar{\zeta} / 2,0, \ldots, 0), \zeta \in \mathbb{T}$.

Since $r_{0}$ is real-analytic on $U_{0} \subset \mathbb{R}^{2 n}$, it extends in a natural way to a holomorphic function in a neighborhood $U_{0}^{\mathbb{C}} \subset \mathbb{C}^{2 n}$ of $U_{0}$. Without loss of
generality we may assume that $r_{0}$ is bounded on $U_{0}^{\mathbb{C}}$. Set

$$
X_{0}=X_{0}\left(U_{0}, U_{0}^{\mathbb{C}}\right):=\left\{r \in \mathcal{O}\left(U_{0}^{\mathbb{C}}\right): r\left(U_{0}\right) \subset \mathbb{R} \text { and } r \text { is bounded }\right\},
$$

which equipped with the sup-norm is a (real) Banach space.
Remark 5.2.1. Lempert considered the case when $U_{0}$ is a neighborhood of the boundary of a bounded domain $D_{0}$ with real-analytic boundary. We shall need more general results to prove the 'localization property'.
5.3. General lemmas. We keep the notation from Subsection 5.2 and assume situation ( $\dagger$ ).

Let us introduce some additional objects we shall be dealing with and let us prove more general lemmas (their generality will be useful in the next section).

Consider the Sobolev space $W^{2,2}\left(\mathbb{T}, \mathbb{C}^{m}\right)$ of functions $f: \mathbb{T} \rightarrow \mathbb{C}^{m}$ whose first two derivatives (in the sense of distributions) are in $L^{2}(\mathbb{T})$. The $W^{2,2_{-}}$ norm is denoted by $\|\cdot\|_{W}$. For the basic properties of $W^{2,2}\left(\mathbb{T}, \mathbb{C}^{m}\right)$ see Appendix.

Put
$B:=\left\{f \in W^{2,2}\left(\mathbb{T}, \mathbb{C}^{n}\right): f\right.$ extends holomorphically on $\mathbb{D}$ and $\left.f(0)=0\right\}$,

$$
\begin{aligned}
B_{0} & :=\left\{f \in B: f(\mathbb{T}) \subset U_{0}\right\}, & B^{*}:=\{\bar{f}: f \in B\}, \\
Q & :=\left\{q \in W^{2,2}(\mathbb{T}, \mathbb{C}): q(\mathbb{T}) \subset \mathbb{R}\right\}, & Q_{0}:=\{q \in Q: q(1)=0\} .
\end{aligned}
$$

It is clear that $B, B^{*}, Q$ and $Q_{0}$ equipped with the norm $\|\cdot\|_{W}$ are (real) Banach spaces. Note that $B_{0}$ is an open neighborhood of $f_{0}$. In what follows, we identify $f \in B$ with its unique holomorphic extension on $\mathbb{D}$.

Let us define the projection

$$
\pi: W^{2,2}\left(\mathbb{T}, \mathbb{C}^{n}\right) \ni f=\sum_{k=-\infty}^{\infty} a_{k} \zeta^{k} \mapsto \sum_{k=-\infty}^{-1} a_{k} \zeta^{k} \in B^{*}
$$

Note that $f \in W^{2,2}\left(\mathbb{T}, \mathbb{C}^{n}\right)$ extends holomorphically on $\mathbb{D}$ if and only if $\pi(f)=0$ (and the extension is $\mathcal{C}^{1 / 2}$ on $\mathbb{T}$ ). Actually, it suffices to observe that $g(\zeta):=\sum_{k=-\infty}^{-1} a_{k} \zeta^{k}, \zeta \in \mathbb{T}$, extends holomorphically on $\mathbb{D}$ if and only if $a_{k}=0$ for $k<0$. This follows immediately from the fact that the mapping $\mathbb{T} \ni \zeta \mapsto g(\bar{\zeta}) \in \mathbb{C}^{n}$ extends holomorphically on $\mathbb{D}$.

Consider the mapping $\Xi: X_{0} \times \mathbb{C}^{n} \times B_{0} \times Q_{0} \times \mathbb{R} \rightarrow Q \times B^{*} \times \mathbb{C}^{n}$ defined by

$$
\Xi(r, v, f, q, \lambda):=\left(r \circ f, \pi\left(\zeta(1+q)\left(r_{z} \circ f\right)\right), f^{\prime}(0)-\lambda v\right),
$$

where $\zeta$ is treated as the identity function on $\mathbb{T}$. We have the following
Lemma 5.3.1. There exist a neighborhood $V_{0}$ of $\left(r_{0}, f_{0}^{\prime}(0)\right)$ in $X_{0} \times \mathbb{C}^{n}$ and a real-analytic mapping $\Upsilon: V_{0} \rightarrow B_{0} \times Q_{0} \times \mathbb{R}$ such that for any $(r, v) \in V_{0}$ we have $\Xi(r, v, \Upsilon(r, v))=0$.

Let $\widetilde{\Xi}: X_{0} \times \mathbb{C}^{n} \times B_{0} \times Q_{0} \times(0,1) \rightarrow Q \times B^{*} \times \mathbb{C}^{n}$ be defined as

$$
\widetilde{\Xi}(r, w, f, q, \xi):=\left(r \circ f, \pi\left(\zeta(1+q)\left(r_{z} \circ f\right)\right), f(\xi)-w\right) .
$$

Analogously we have
Lemma 5.3.2. Let $\xi_{0} \in(0,1)$. Then there exist a neighborhood $W_{0}$ of $\left(r_{0}, f_{0}\left(\xi_{0}\right)\right)$ in $X_{0} \times D_{0}$ and a real-analytic mapping $\widetilde{\Upsilon}: W_{0} \rightarrow B_{0} \times Q_{0} \times(0,1)$ such that for any $(r, w) \in W_{0}$ we have $\widetilde{\Xi}(r, w, \widetilde{\Upsilon}(r, w))=0$.

Proof of Lemmas 5.3.1 and 5.3.2. We will prove the first lemma. Then we will see that the proof of the second one reduces to that proof.

We claim that $\Xi$ is real-analytic. The only problem is to show that the mapping

$$
T: X_{0} \times B_{0} \ni(r, f) \mapsto r \circ f \in Q
$$

is real-analytic (the real-analyticity of the mapping $X_{0} \times B_{0} \ni(r, f) \mapsto$ $r_{z} \circ f \in W^{2,2}\left(\mathbb{T}, \mathbb{C}^{n}\right)$ follows from this claim).

Fix $r \in X_{0}, f \in B_{0}$ and take $\varepsilon>0$ such that a $2 n$-dimensional polydisc $P_{2 n}(f(\zeta), \varepsilon)$ is contained in $U_{0}^{\mathbb{C}}$ for any $\zeta \in \mathbb{T}$. Then any function $\widetilde{r} \in X_{0}$ is holomorphic in $U_{0}^{\mathbb{C}}$, so it expands in a holomorphic power series convergent in $P_{2 n}(f(\zeta), \varepsilon)$. Losing no generality we may assume that $n$-dimensional polydiscs $P_{n}(f(\zeta), \varepsilon), \zeta \in \mathbb{T}$, satisfy $P_{n}(f(\zeta), \varepsilon) \subset U_{0}$. This gives an expansion of the function $\widetilde{r}$ at any point $f(\zeta), \zeta \in \mathbb{T}$, in a series

$$
\sum_{\alpha \in \mathbb{N}_{0}^{2 n}} \frac{1}{\alpha!} \frac{\partial^{|\alpha|} \widetilde{r}}{\partial x^{\alpha}}(f(\zeta)) x^{\alpha}
$$

convergent to $\widetilde{r}(f(\zeta)+x)$, provided that $x=\left(x_{1}, \ldots, x_{2 n}\right) \in P_{n}(0, \varepsilon)$ (where $\mathbb{N}_{0}:=\mathbb{N} \cup\{0\}$ and $\left.|\alpha|:=\alpha_{1}+\cdots+\alpha_{2 n}\right)$. Hence

$$
\begin{equation*}
T(r+\varrho, f+h)=\sum_{\alpha \in \mathbb{N}_{0}^{2 n}} \frac{1}{\alpha!}\left(\frac{\partial^{|\alpha|} r}{\partial x^{\alpha}} \circ f\right) h^{\alpha}+\sum_{\alpha \in \mathbb{N}_{0}^{2 n}} \frac{1}{\alpha!}\left(\frac{\partial^{|\alpha|} \varrho}{\partial x^{\alpha}} \circ f\right) h^{\alpha} \tag{22}
\end{equation*}
$$

pointwise for $\varrho \in X_{0}$ and $h \in W^{2,2}\left(\mathbb{T}, \mathbb{C}^{n}\right)$ with $\|h\|_{\text {sup }}<\varepsilon$.
Put $P:=\bigcup_{\zeta \in \mathbb{T}} P_{2 n}(f(\zeta), \varepsilon)$ and for $\widetilde{r} \in X_{0}$ put $\|\widetilde{r}\|_{P}:=\sup _{P}|\widetilde{r}|$. Let $\widetilde{r}$ be equal to $r$ or to $\varrho$, where $\varrho$ lies in a neighborhood of 0 in $X_{0}$. The Cauchy inequalities give

$$
\begin{equation*}
\left|\frac{\partial^{|\alpha|} \widetilde{r}}{\partial x^{\alpha}}(f(\zeta))\right| \leq \frac{\alpha!\|\widetilde{r}\|_{P}}{\varepsilon^{|\alpha|}}, \quad \zeta \in \mathbb{T} . \tag{23}
\end{equation*}
$$

Therefore,

$$
\left\|\frac{\partial^{|\alpha|} \widetilde{r}}{\partial x^{\alpha}} \circ f\right\|_{W} \leq C_{1} \frac{\alpha!\mid\|\widetilde{r}\|_{P}}{\varepsilon^{|\alpha|}}
$$

for some $C_{1}>0$.

There is $C_{2}>0$ such that

$$
\left\|g h^{\alpha}\right\|_{W} \leq C_{2}^{|\alpha|+1}\|g\|_{W}\left\|h_{1}\right\|_{W}^{\alpha_{1}} \cdot \ldots \cdot\left\|h_{2 n}\right\|_{W}^{\alpha_{2 n}}
$$

for $g \in W^{2,2}(\mathbb{T}, \mathbb{C}), h \in W^{2,2}\left(\mathbb{T}, \mathbb{C}^{n}\right), \alpha \in \mathbb{N}_{0}^{2 n}$ (see Appendix for a proof of this fact). Using the above inequalities we infer that

$$
\sum_{\alpha \in \mathbb{N}_{0}^{2 n}}\left\|\frac{1}{\alpha!}\left(\frac{\partial^{|\alpha|} \widetilde{r}}{\partial x^{\alpha}} \circ f\right) h^{\alpha}\right\|_{W}
$$

is convergent if $\|h\|_{W}$ is small enough. Therefore, the series in 22) are absolutely convergent in the norm $\|\cdot\|_{W}$, whence $T$ is real-analytic.

To show the existence of $V_{0}$ and $\Upsilon$ we will make use of the Implicit Function Theorem. More precisely, we shall show that the partial derivative

$$
\Xi_{(f, q, \lambda)}\left(r_{0}, f_{0}^{\prime}(0), f_{0}, 0,1\right): B \times Q_{0} \times \mathbb{R} \rightarrow Q \times B^{*} \times \mathbb{C}^{n}
$$

is an isomorphism. Observe that for any $(\tilde{f}, \widetilde{q}, \tilde{\lambda}) \in B \times Q_{0} \times \mathbb{R}$ the following equality holds:

$$
\begin{array}{r}
\Xi_{(f, q, \lambda)}\left(r_{0}, f_{0}^{\prime}(0), f_{0}, 0,1\right)(\widetilde{f}, \widetilde{q}, \widetilde{\lambda})=\left.\frac{d}{d t} \Xi\left(r_{0}, f_{0}^{\prime}(0), f_{0}+t \widetilde{f}, t \widetilde{q}, 1+t \widetilde{\lambda}\right)\right|_{t=0} \\
=\left(\left(r_{0 z} \circ f_{0}\right) \widetilde{f}+\left(r_{0 \bar{z}} \circ f_{0}\right) \overline{\widetilde{f}}, \pi\left(\zeta \widetilde{q} r_{0 z} \circ f_{0}+\zeta\left(r_{0 z z} \circ f_{0}\right) \widetilde{f}+\zeta\left(r_{0 z \bar{z}} \circ f_{0}\right) \overline{\widetilde{f}}\right)\right. \\
\left.\widetilde{f^{\prime}}(0)-\widetilde{\lambda} f_{0}^{\prime}(0)\right)
\end{array}
$$

where we treat $r_{0 z}, r_{0 \bar{z}}$ as row vectors, $\widetilde{f}, \overline{\widetilde{f}}$ as column vectors and $r_{0 z z}=$ $\left[\frac{\partial^{2} r_{0}}{\partial z_{j} \partial z_{k}}\right]_{j, k=1}^{n}, r_{0 z \bar{z}}=\left[\frac{\partial^{2} r_{0}}{\partial z_{j} \partial \bar{z}_{k}}\right]_{j, k=1}^{n}$ as $n \times n$ matrices.

By the Bounded Inverse Theorem, it suffices to show that the derivative $\Xi_{(f, q, \lambda)}\left(r_{0}, f_{0}^{\prime}(0), f_{0}, 0,1\right)$ is bijective, i.e. for any $(\eta, \varphi, v) \in Q \times B^{*} \times \mathbb{C}^{n}$ there exists exactly one $(\widetilde{f}, \widetilde{q}, \widetilde{\lambda}) \in B \times Q_{0} \times \mathbb{R}$ satisfying

$$
\begin{align*}
\left(r_{0 z} \circ f_{0}\right) \widetilde{f}+\left(r_{0 \bar{z}} \circ f_{0}\right) \overline{\widetilde{f}} & =\eta  \tag{24}\\
\pi\left(\zeta \widetilde{q} r_{0 z} \circ f_{0}+\zeta\left(r_{0 z z} \circ f_{0}\right) \widetilde{f}+\zeta\left(r_{0 z \bar{z}} \circ f_{0}\right) \overline{\widetilde{f}}\right) & =\varphi  \tag{25}\\
\widetilde{f^{\prime}}(0)-\widetilde{\lambda} f_{0}^{\prime}(0) & =v \tag{26}
\end{align*}
$$

First, we show that $\widetilde{\lambda}$ and $\widetilde{f}_{1}$ are uniquely determined. Observe that, in view of the assumptions, (24) is just

$$
\frac{1}{2} \bar{\zeta} \widetilde{f}_{1}+\frac{1}{2} \zeta \widetilde{\widetilde{f}}_{1}=\eta
$$

or equivalently

$$
\begin{equation*}
\operatorname{Re}\left(\tilde{f}_{1} / \zeta\right)=\eta \quad(\text { on } \mathbb{T}) \tag{27}
\end{equation*}
$$

Note that 27 uniquely determines $\tilde{f}_{1} / \zeta \in W^{2,2}(\mathbb{T}, \mathbb{C}) \cap \mathcal{O}(\mathbb{D}) \cap \mathcal{C}(\overline{\mathbb{D}})$ up to an imaginary additive constant, which may be computed using (26). Actually, $\eta=\operatorname{Re} G$ on $\mathbb{T}$ for some function $G \in W^{2,2}(\mathbb{T}, \mathbb{C}) \cap \mathcal{O}(\mathbb{D}) \cap \mathcal{C}(\overline{\mathbb{D}})$. To see this,
let us expand $\eta(\zeta)=\sum_{k=-\infty}^{\infty} a_{k} \zeta^{k}, \zeta \in \mathbb{T}$. From the equality $\eta(\zeta)=\overline{\eta(\zeta)}$, $\zeta \in \mathbb{T}$, we get

$$
\begin{equation*}
\sum_{k=-\infty}^{\infty} a_{k} \zeta^{k}=\sum_{k=-\infty}^{\infty} \bar{a}_{k} \zeta^{-k}=\sum_{k=-\infty}^{\infty} \bar{a}_{-k} \zeta^{k}, \quad \zeta \in \mathbb{T} \tag{28}
\end{equation*}
$$

so $a_{-k}=\bar{a}_{k}, k \in \mathbb{Z}$. Hence

$$
\eta(\zeta)=a_{0}+\sum_{k=1}^{\infty} 2 \operatorname{Re}\left(a_{k} \zeta^{k}\right)=\operatorname{Re}\left(a_{0}+2 \sum_{k=1}^{\infty} a_{k} \zeta^{k}\right), \quad \zeta \in \mathbb{T}
$$

Set

$$
G(\zeta):=a_{0}+2 \sum_{k=1}^{\infty} a_{k} \zeta^{k}, \quad \zeta \in \mathbb{D}
$$

This series is convergent for $\zeta \in \mathbb{D}$, so $G \in \mathcal{O}(\mathbb{D})$. Further, the function $G$ extends continuously on $\overline{\mathbb{D}}$ (to the function denoted by the same letter) and the extension lies in $W^{2,2}(\mathbb{T}, \mathbb{C})$. Clearly, $\eta=\operatorname{Re} G$ on $\mathbb{T}$.

We are searching for $C \in \mathbb{R}$ such that the functions $\widetilde{f}_{1}:=\zeta(G+i C)$ and $\theta:=\operatorname{Im}\left(\widetilde{f}_{1} / \zeta\right)$ satisfy

$$
\eta(0)+i \theta(0)={\widetilde{f_{1}^{\prime}}}_{1}^{\prime}(0)
$$

and

$$
\eta(0)+i \theta(0)-\widetilde{\lambda} \operatorname{Re} f_{01}^{\prime}(0)-i \widetilde{\lambda} \operatorname{Im} f_{01}^{\prime}(0)=\operatorname{Re} v_{1}+i \operatorname{Im} v_{1}
$$

But

$$
\eta(0)-\widetilde{\lambda} \operatorname{Re} f_{01}^{\prime}(0)=\operatorname{Re} v_{1}
$$

which yields $\widetilde{\lambda}$ and then $\theta(0)$, and consequently the number $C$. Having $\widetilde{\lambda}$ and once again using 26 , we find uniquely determined $\widetilde{f}_{2}^{\prime}(0), \ldots, \widetilde{f}_{n}^{\prime}(0)$.

Therefore, the equations (24) and (26) are satisfied by uniquely determined $\widetilde{f}_{1}, \widetilde{\lambda}$ and $\widetilde{f}_{2}^{\prime}(0), \ldots, \widetilde{f}_{n}^{\prime}(0)$.

Consider 25 ), which is a system of $n$ equations with unknowns $\widetilde{q}, \widetilde{f}_{2}, \ldots, \tilde{f}_{n}$. Observe that $\widetilde{q}$ appears only in the first of the equations and the remaining $n-1$ equations mean exactly that the mapping

$$
\begin{equation*}
\zeta\left(r_{0 \widehat{z} \widehat{z}} \circ f_{0}\right) \widehat{\widetilde{f}}+\zeta\left(r_{0 \widehat{z} \widehat{\bar{z}}} \circ f_{0}\right) \widehat{\overline{\tilde{f}}}-\psi \tag{29}
\end{equation*}
$$

extends holomorphically on $\mathbb{D}$; here $\widehat{a}:=\left(a_{2}, \ldots, a_{n}\right)$ and $\psi \in W^{2,2}\left(\mathbb{T}, \mathbb{C}^{n-1}\right)$ may be obtained from $\varphi$ and $\widetilde{f}_{1}$. To see this, write 25 in the form

$$
\pi\left(F_{1}+\zeta F_{2}+\zeta F_{3}\right)=\left(\varphi_{1}, \ldots, \varphi_{n}\right)
$$

where

$$
\begin{aligned}
& F_{1}:=(\widetilde{q}, 0, \ldots, 0) \\
& F_{2}:=\left(A_{j}\right)_{j=1}^{n}, \quad A_{j}:=\sum_{k=1}^{n}\left(r_{0 z_{j} z_{k}} \circ f_{0}\right) \widetilde{f}_{k}, \\
& F_{3}:=\left(B_{j}\right)_{j=1}^{n}, \quad B_{j}:=\sum_{k=1}^{n}\left(r_{0 z_{j} \bar{z}_{k}} \circ f_{0}\right) \overline{\widetilde{f}}_{k} .
\end{aligned}
$$

It follows that

$$
\widetilde{q}+\zeta A_{1}+\zeta B_{1}-\varphi_{1} \quad \text { and } \quad \zeta A_{j}+\zeta B_{j}-\varphi_{j}, \quad j=2, \ldots, n
$$

should extend holomorphically on $\mathbb{D}$ and

$$
\psi:=\left(\varphi_{j}-\zeta\left(r_{0 z_{j} z_{1}} \circ f_{0}\right) \widetilde{f}_{1}-\zeta\left(r_{0 z_{j} \bar{z}_{1}} \circ f_{0}\right) \overline{\widetilde{f}}_{1}\right)_{j=2}^{n}
$$

Put

$$
g(\zeta):=\widehat{\widetilde{f}}(\zeta) / \zeta, \quad \alpha(\zeta):=\zeta^{2} r_{0 \widetilde{z} \bar{z}}\left(f_{0}(\zeta)\right), \quad \beta(\zeta):=r_{0 \widehat{z} \widehat{\bar{z}}}\left(f_{0}(\zeta)\right)
$$

Observe that $\alpha(\zeta), \beta(\zeta)$ are $(n-1) \times(n-1)$ matrices depending realanalytically on $\zeta$, and $g(\zeta)$ is a column vector in $\mathbb{C}^{n-1}$. This allows us to reduce (29) to the following problem: find a unique $g \in W^{2,2}\left(\mathbb{T}, \mathbb{C}^{n-1}\right) \cap \mathcal{O}(\mathbb{D}) \cap \mathcal{C}(\overline{\mathbb{D}})$ such that

$$
\begin{equation*}
\alpha g+\beta \bar{g}-\psi \text { extends holomorphically on } \mathbb{D} \text { and } g(0)=\widehat{\widetilde{f^{\prime}}}(0) \tag{30}
\end{equation*}
$$

The fact that every $f_{0}(\zeta)$ is a point of strong linear convexity of $D_{0}$ may be written as

$$
\begin{equation*}
\left|X^{T} \alpha(\zeta) X\right|<X^{T} \beta(\zeta) \bar{X}, \quad \zeta \in \mathbb{T}, X \in\left(\mathbb{C}^{n-1}\right)_{*} \tag{31}
\end{equation*}
$$

Note that $\beta(\zeta)$ is self-adjoint and strictly positive, hence using Proposition 8.4.1 we get a mapping $H \in \mathcal{O}\left(\overline{\mathbb{D}}, \mathbb{C}^{(n-1) \times(n-1)}\right)$ such that $\operatorname{det} H \neq 0$ on $\overline{\mathbb{D}}$ and $H H^{*}=\beta$ on $\mathbb{T}$. With this notation, (30) is equivalent to

$$
\begin{equation*}
H^{-1} \alpha g+H^{*} \bar{g}-H^{-1} \psi \text { extends holomorphically on } \mathbb{D} \tag{32}
\end{equation*}
$$

or, if we denote $h:=H^{T} g, \gamma:=H^{-1} \alpha\left(H^{T}\right)^{-1}$, to

$$
\begin{equation*}
\gamma h+\bar{h}-H^{-1} \psi \text { extends holomorphically on } \mathbb{D} . \tag{33}
\end{equation*}
$$

For any $\zeta \in \mathbb{T}$ the operator norm of the symmetric matrix $\gamma(\zeta)$ is uniformly less than 1. In fact, from 31), for any $X \in \mathbb{C}^{n-1}$ with $|X|=1$,

$$
\begin{aligned}
\left|X^{T} \gamma(\zeta) X\right| & =\left|X^{T} H(\zeta)^{-1} \alpha(\zeta)\left(H(\zeta)^{T}\right)^{-1} X\right| \\
& <X^{T} H(\zeta)^{-1} \beta(\zeta) \overline{\left(H(\zeta)^{T}\right)^{-1} X} \\
& =X^{T} H(\zeta)^{-1} H(\zeta) H(\zeta)^{*} \overline{\left(H(\zeta)^{T}\right)^{-1}} \bar{X}=|X|^{2}=1
\end{aligned}
$$

so, by a compactness argument, $\left|X^{T} \gamma(\zeta) X\right| \leq 1-\widetilde{\varepsilon}$ for some $\widetilde{\varepsilon}>0$ independent of $\zeta$ and $X$. Thus $\|\gamma(\zeta)\| \leq 1-\widetilde{\varepsilon}$ by Proposition 8.4.2.

We have to prove that there is a unique solution $h \in W^{2,2}\left(\mathbb{T}, \mathbb{C}^{n-1}\right) \cap$ $\mathcal{O}(\mathbb{D}) \cap \mathcal{C}(\overline{\mathbb{D}})$ of 33 such that $h(0)=a$ with a given $a \in \mathbb{C}^{n-1}$.

Define the operator

$$
P: W^{2,2}\left(\mathbb{T}, \mathbb{C}^{n-1}\right) \ni \sum_{k=-\infty}^{\infty} a_{k} \zeta^{k} \mapsto \overline{\sum_{k=-\infty}^{-1} a_{k} \zeta^{k}} \in W^{2,2}\left(\mathbb{T}, \mathbb{C}^{n-1}\right)
$$

where $a_{k} \in \mathbb{C}^{n-1}, k \in \mathbb{Z}$. We will show that a mapping $h \in W^{2,2}\left(\mathbb{T}, \mathbb{C}^{n-1}\right) \cap$ $\mathcal{O}(\mathbb{D}) \cap \mathcal{C}(\overline{\mathbb{D}})$ satisfies 33 and $h(0)=a$ if and only if it is a fixed point of the mapping

$$
K: W^{2,2}\left(\mathbb{T}, \mathbb{C}^{n-1}\right) \ni h \mapsto P\left(H^{-1} \psi-\gamma h\right)+a \in W^{2,2}\left(\mathbb{T}, \mathbb{C}^{n-1}\right)
$$

Indeed, take $h \in W^{2,2}\left(\mathbb{T}, \mathbb{C}^{n-1}\right) \cap \mathcal{O}(\mathbb{D}) \cap \mathcal{C}(\overline{\mathbb{D}})$ and suppose that $h(0)=a$ and $\gamma h+\bar{h}-H^{-1} \psi$ extends holomorphically on $\mathbb{D}$. Then

$$
\begin{array}{ll}
h=a+\sum_{k=1}^{\infty} a_{k} \zeta^{k}, & \bar{h}=\bar{a}+\sum_{k=1}^{\infty} \bar{a}_{k} \zeta^{-k}=\sum_{k=-\infty}^{-1} \bar{a}_{-k} \zeta^{k}+\bar{a} \\
P(h)=0, & P(\bar{h})=\sum_{k=1}^{\infty} a_{k} \zeta^{k}=h-a
\end{array}
$$

and

$$
P\left(\gamma h+\bar{h}-H^{-1} \psi\right)=0
$$

which implies

$$
P\left(H^{-1} \psi-\gamma h\right)=h-a
$$

and finally $K(h)=h$. Conversely, suppose that $K(h)=h$. Then

$$
P\left(H^{-1} \psi-\gamma h\right)=h-a=\sum_{k=1}^{\infty} a_{k} \zeta^{k}+a_{1}-a, \quad P(h)=0
$$

and

$$
P(\bar{h})=\sum_{k=1}^{\infty} a_{k} \zeta^{k}=h-a_{1}
$$

from which follows that

$$
P\left(\gamma h+\bar{h}-H^{-1} \psi\right)=P(\bar{h})-P\left(H^{-1} \psi-\gamma h\right)=a-a_{1}
$$

and

$$
P\left(\gamma h+\bar{h}-H^{-1} \psi\right)=0 \quad \text { iff } \quad a=a_{1}
$$

Observe that $h(0)=K(h)(0)=P\left(H^{-1} \psi-\gamma h\right)(0)+a=a$.
We shall make use of the Banach Fixed Point Theorem. To do this, consider $W^{2,2}\left(\mathbb{T}, \mathbb{C}^{n-1}\right)$ equipped with the following norm:

$$
\|h\|_{\varepsilon}:=\|h\|_{L}+\varepsilon\left\|h^{\prime}\right\|_{L}+\varepsilon^{2}\left\|h^{\prime \prime}\right\|_{L},
$$

where $\varepsilon>0$ and $\|\cdot\|_{L}$ is the $L^{2}$-norm (it is a Banach space). We will prove that $K$ is a contraction with respect to the norm $\|\cdot\|_{\varepsilon}$ for sufficiently small $\varepsilon$. Indeed, there is $\widetilde{\varepsilon}>0$ such that, for any $h_{1}, h_{2} \in W^{2,2}\left(\mathbb{T}, \mathbb{C}^{n-1}\right)$,

$$
\begin{align*}
\left\|K\left(h_{1}\right)-K\left(h_{2}\right)\right\|_{L} & =\left\|P\left(\gamma\left(h_{2}-h_{1}\right)\right)\right\|_{L}  \tag{34}\\
& \leq\left\|\gamma\left(h_{2}-h_{1}\right)\right\|_{L} \leq(1-\widetilde{\varepsilon})\left\|h_{2}-h_{1}\right\|_{L} .
\end{align*}
$$

Moreover,

$$
\begin{align*}
\| K\left(h_{1}\right)^{\prime}- & K\left(h_{2}\right)^{\prime} \|_{L} \tag{35}
\end{align*} \quad\left\|P\left(\gamma h_{2}\right)^{\prime}-P\left(\gamma h_{1}\right)^{\prime}\right\|_{L} .
$$

Furthermore,

$$
\begin{align*}
\| K\left(h_{1}\right)^{\prime \prime}- & K\left(h_{2}\right)^{\prime \prime} \|_{L}  \tag{36}\\
& \leq\left\|\gamma^{\prime \prime}\left(h_{2}-h_{1}\right)\right\|_{L}+2\left\|\gamma^{\prime}\left(h_{2}^{\prime}-h_{1}^{\prime}\right)\right\|_{L}+\left\|\gamma\left(h_{1}^{\prime \prime}-h_{2}^{\prime \prime}\right)\right\|_{L} .
\end{align*}
$$

Using the finiteness of $\left\|\gamma^{\prime}\right\|,\left\|\gamma^{\prime \prime}\right\|$ and putting (34), (35), (36) together, we see that there exists $\varepsilon>0$ such that $K$ is a contraction with respect to $\|\cdot\|_{\varepsilon}$.

We have found $\tilde{f}$ and $\tilde{\lambda}$ satisfying (24), (26), and the last $n-1$ equations from (25) are satisfied as well.

It remains to show that there exists a unique $\widetilde{q} \in Q_{0}$ such that $\widetilde{q}+\zeta A_{1}+$ $\zeta B_{1}-\varphi_{1}$ extends holomorphically on $\mathbb{D}$.

Comparing the coefficients as in (28), we see that if

$$
\pi\left(\zeta A_{1}+\zeta B_{1}-\varphi_{1}\right)=\sum_{k=-\infty}^{-1} a_{k} \zeta^{k}
$$

then $\widetilde{q}$ has to be taken as

$$
-\sum_{k=-\infty}^{-1} a_{k} \zeta^{k}-\sum_{k=0}^{\infty} b_{k} \zeta^{k}
$$

with $b_{k}:=\bar{a}_{-k}$ for $k \geq 1$ and $b_{0} \in \mathbb{R}$ uniquely determined by $\widetilde{q}(1)=0$.
Let us show that Lemma 5.3.2 follows from the proof of 5.3.1. Since $\widetilde{\Xi}$ is real-analytic, it suffices to prove that the derivative

$$
\widetilde{\Xi}_{(f, q, \xi)}\left(r_{0}, f_{0}\left(\xi_{0}\right), f_{0}, 0, \xi_{0}\right): B \times Q_{0} \times \mathbb{R} \rightarrow Q \times B^{*} \times \mathbb{C}^{n}
$$

is invertible. For $(\widetilde{f}, \widetilde{q}, \widetilde{\xi}) \in B \times Q_{0} \times \mathbb{R}$ we get

$$
\begin{array}{r}
\widetilde{\Xi}_{(f, q, \xi)}\left(r_{0}, f_{0}\left(\xi_{0}\right), f_{0}, 0, \xi_{0}\right)(\widetilde{f}, \widetilde{q}, \widetilde{\xi})=\left.\frac{d}{d t} \widetilde{\Xi}\left(r_{0}, f_{0}\left(\xi_{0}\right), f_{0}+t \widetilde{f}, t \widetilde{q}, \xi_{0}+t \widetilde{\xi}\right)\right|_{t=0} \\
=\left(\left(r_{0 z} \circ f_{0}\right) \widetilde{f}+\left(r_{0 \bar{z}} \circ f_{0}\right) \overline{\widetilde{f}}, \pi\left(\zeta \widetilde{q} r_{0 z} \circ f_{0}+\zeta\left(r_{0 z z} \circ f_{0}\right) \widetilde{f}+\zeta\left(r_{0 z \bar{z}} \circ f_{0}\right) \tilde{\tilde{f}}\right),\right. \\
\left.\widetilde{f}\left(\xi_{0}\right)+\widetilde{\xi} f_{0}^{\prime}\left(\xi_{0}\right)\right) .
\end{array}
$$

We have to show that for any $(\eta, \varphi, w) \in Q \times B^{*} \times \mathbb{C}^{n}$ there exists exactly one $(\widetilde{f}, \widetilde{q}, \widetilde{\xi}) \in B \times Q_{0} \times \mathbb{R}$ satisfying

$$
\begin{align*}
\left(r_{0 z} \circ f_{0}\right) \widetilde{f}+\left(r_{0 \bar{z}} \circ f_{0}\right) \overline{\widetilde{f}} & =\eta  \tag{37}\\
\pi\left(\zeta \widetilde{q} r_{0 z} \circ f_{0}+\zeta\left(r_{0 z z} \circ f_{0}\right) \widetilde{f}+\zeta\left(r_{0 z \bar{z}} \circ f_{0}\right) \overline{\widetilde{f}}\right) & =\varphi  \tag{38}\\
\widetilde{f}\left(\xi_{0}\right)+\widetilde{\xi} f_{0}^{\prime}\left(\xi_{0}\right) & =w \tag{39}
\end{align*}
$$

The equation (37) turns out to be

$$
\begin{equation*}
\operatorname{Re}\left(\tilde{f}_{1} / \zeta\right)=\eta \quad(\text { on } \mathbb{T}) \tag{40}
\end{equation*}
$$

The equation above uniquely determines $\widetilde{f}_{1} / \zeta \in W^{2,2}(\mathbb{T}, \mathbb{C}) \cap \mathcal{O}(\mathbb{D}) \cap \mathcal{C}(\overline{\mathbb{D}})$ up to an imaginary additive constant, which may be computed using (39). Indeed, there exists $G \in W^{2,2}(\mathbb{T}, \mathbb{C}) \cap \mathcal{O}(\mathbb{D}) \cap \mathcal{C}(\overline{\mathbb{D}})$ such that $\eta=\operatorname{Re} G$ on $\mathbb{T}$. We are searching for $C \in \mathbb{R}$ such that the functions $\widetilde{f}_{1}:=\zeta(G+i C)$ and $\theta:=\operatorname{Im}\left(\widetilde{f}_{1} / \zeta\right)$ satisfy

$$
\xi_{0} \eta\left(\xi_{0}\right)+i \xi_{0} \theta\left(\xi_{0}\right)=\widetilde{f}_{1}\left(\xi_{0}\right)
$$

and

$$
\xi_{0}\left(\eta\left(\xi_{0}\right)+i \theta\left(\xi_{0}\right)\right)+\widetilde{\xi} \operatorname{Re} f_{01}^{\prime}\left(\xi_{0}\right)+i \widetilde{\xi} \operatorname{Im} f_{01}^{\prime}\left(\xi_{0}\right)=\operatorname{Re} w_{1}+i \operatorname{Im} w_{1}
$$

But

$$
\xi_{0} \eta\left(\xi_{0}\right)+\widetilde{\xi} \operatorname{Re} f_{01}^{\prime}\left(\xi_{0}\right)=\operatorname{Re} w_{1}
$$

which yields $\widetilde{\xi}$ and then $\theta\left(\xi_{0}\right)$, and consequently the number $C$. Having $\widetilde{\xi}$ and once again using 39, we find uniquely determined $\widetilde{f}_{2}\left(\xi_{0}\right), \ldots, \widetilde{f}_{n}\left(\xi_{0}\right)$.

Therefore, the equations (37) and (39) are satisfied by uniquely determined $\widetilde{f}_{1}, \widetilde{\xi}$ and $\widetilde{f}_{2}\left(\xi_{0}\right), \ldots, \widetilde{f}_{n}\left(\xi_{0}\right)$.

In the remaining part of the proof we change the second condition of 30 to

$$
g\left(\xi_{0}\right)=\widehat{\widetilde{f}}\left(\xi_{0}\right) / \xi_{0}
$$

and we have to prove that there is a unique solution $h \in W^{2,2}\left(\mathbb{T}, \mathbb{C}^{n-1}\right) \cap$ $\mathcal{O}(\mathbb{D}) \cap \mathcal{C}(\overline{\mathbb{D}})$ of 33 such that $h\left(\xi_{0}\right)=a$ with a given $a \in \mathbb{C}^{n-1}$. Let $\tau$ be an automorphism of $\mathbb{D}$ (so it extends holomorphically near $\overline{\mathbb{D}}$ ), which maps 0 to $\xi_{0}$, i.e.

$$
\tau(\xi):=\frac{\xi_{0}-\xi}{1-\bar{\xi}_{0} \xi}, \quad \xi \in \mathbb{D}
$$

Let the mappings $P, K$ be as before. Then $h \in W^{2,2}\left(\mathbb{T}, \mathbb{C}^{n-1}\right) \cap \mathcal{O}(\mathbb{D}) \cap \mathcal{C}(\overline{\mathbb{D}})$ satisfies 33 and $h\left(\xi_{0}\right)=a$ if and only if $h \circ \tau \in W^{2,2}\left(\mathbb{T}, \mathbb{C}^{n-1}\right) \cap \mathcal{O}(\mathbb{D}) \cap \mathcal{C}(\overline{\mathbb{D}})$ satisfies (33) and $(h \circ \tau)(0)=a$. We already know that there is exactly one $\widetilde{h} \in W^{2,2}\left(\mathbb{T}, \mathbb{C}^{n-1}\right) \cap \mathcal{O}(\mathbb{D}) \cap \mathcal{C}(\overline{\mathbb{D}})$ satisfying 33 and $\widetilde{h}(0)=a$. Setting $h:=\widetilde{h} \circ \tau^{-1}$, we get the claim.
5.4. Topology in a class of domains with real-analytic boundaries. We introduce a concept of a domain being close to another domain. Let $D_{0} \subset \mathbb{C}^{n}$ be a bounded domain with real-analytic boundary. Then there exist a neighborhood $U_{0}$ of $\partial D_{0}$ and a real-analytic defining function $r_{0}: U_{0} \rightarrow \mathbb{R}$ such that $\nabla r_{0}$ does not vanish in $U_{0}$ and

$$
D_{0} \cap U_{0}=\left\{z \in U_{0}: r_{0}(z)<0\right\} .
$$

Definition 5.4.1. We say that domains $D$ tend to $D_{0}$ if one can choose their defining functions $r \in X_{0}$ so that they tend to $r_{0}$ in $X_{0}$.

Remark 5.4.2. If $r \in X_{0}$ is near $r_{0}$ with respect to the topology on $X_{0}$, then $\left\{z \in U_{0}: r(z)=0\right\}$ is a compact real-analytic hypersurface which bounds a bounded domain. We denote it by $D^{r}$.

Moreover, if $D^{r_{0}}$ is strongly linearly convex then so is $D^{r}$ provided that $r$ is near $r_{0}$.

### 5.5. Statement of the main result of this section

Remark 5.5.1. Assume that $D^{r}$ is a strongly linearly convex domain bounded by a real-analytic hypersurface $\left\{z \in U_{0}: r(z)=0\right\}$. Let $\xi \in(0,1)$ and $w \in\left(\mathbb{C}^{n}\right)_{*}$. Then a function $f \in B_{0}$ satisfies the conditions
$f$ is a weak stationary mapping of $D^{r}, \quad f(0)=0, \quad f(\xi)=w$
if and only if there exists $q \in Q_{0}$ such that $q>-1$ and $\widetilde{\Xi}(r, w, f, q, \xi)=0$.
Actually, from $\widetilde{\Xi}(r, w, f, q, \xi)=0$ we deduce immediately that $r \circ f=0$ on $\mathbb{T}, f(\xi)=w$ and $\pi\left(\zeta(1+q)\left(r_{z} \circ f\right)\right)=0$. From the first equality we get $f(\mathbb{T}) \subset \partial D^{r}$. From the last one we infer that the condition ( $3^{\prime}$ ) of Definition 1.11 is satisfied (with $\left.\rho:=(1+q)\left|r_{z} \circ f\right|\right)$. Since $D^{r}$ is strongly linearly convex, $\overline{D^{r}}$ is polynomially convex (use the fact that projections of $\mathbb{C}$-convex domains are $\mathbb{C}$-convex as well, and that $D^{r}$ is smooth). In particular,

$$
f(\overline{\mathbb{D}})=f(\widehat{\mathbb{T}}) \subset \widehat{f(\mathbb{T})} \subset \widehat{\widehat{D^{r}}}=\overline{D^{r}}
$$

where $\widehat{S}:=\left\{z \in \mathbb{C}^{m}:|P(z)| \leq \sup _{S}|P|\right.$ for any $\left.P \in \mathbb{C}\left[z_{1}, \ldots, z_{m}\right]\right\}$ is the polynomial hull of a set $S \subset \mathbb{C}^{m}$.

Note that this implies $f(\mathbb{D}) \subset D^{r}$, by virtue of the fact that $\partial D^{r}$ does not contain non-constant analytic discs (as $D^{r}$ is strongly pseudoconvex).

The opposite implication is clear.
In a similar way we show that for any $v \in\left(\mathbb{C}^{n}\right)_{*}$ and $\lambda>0$, a function $f \in B_{0}$ satisfies the conditions

$$
f \text { is a weak stationary mapping of } D^{r}, \quad f(0)=0, \quad f^{\prime}(0)=\lambda v
$$

if and only if there exists $q \in Q_{0}$ such that $q>-1$ and $\Xi(r, v, f, q, \lambda)=0$.

Proposition 5.5.2. Let $D_{0} \subset \mathbb{C}^{n}$, $n \geq 2$, be a strongly linearly convex domain with real-analytic boundary and let $f_{0}: \mathbb{D} \rightarrow D_{0}$ be an E-mapping.
(1) Let $\xi_{0} \in(0,1)$. Then there exist a neighborhood $W_{0}$ of $\left(r_{0}, f_{0}\left(\xi_{0}\right)\right)$ in $X_{0} \times D_{0}$ and real-analytic mappings

$$
\Lambda: W_{0} \rightarrow \mathcal{C}^{1 / 2}(\overline{\mathbb{D}}), \quad \Omega: W_{0} \rightarrow(0,1)
$$

such that

$$
\Lambda\left(r_{0}, f_{0}\left(\xi_{0}\right)\right)=f_{0}, \quad \Omega\left(r_{0}, f_{0}\left(\xi_{0}\right)\right)=\xi_{0}
$$

and for any $(r, w) \in W_{0}$ the mapping $f:=\Lambda(r, w)$ is an E-mapping of $D^{r}$ satisfying

$$
f(0)=f_{0}(0) \quad \text { and } \quad f(\Omega(r, w))=w
$$

(2) There exist a neighborhood $V_{0}$ of $\left(r_{0}, f_{0}^{\prime}(0)\right)$ in $X_{0} \times \mathbb{C}^{n}$ and a realanalytic mapping

$$
\Gamma: V_{0} \rightarrow \mathcal{C}^{1 / 2}(\overline{\mathbb{D}})
$$

such that

$$
\Gamma\left(r_{0}, f_{0}^{\prime}(0)\right)=f_{0}
$$

and for any $(r, v) \in V_{0}$ the mapping $f:=\Gamma(r, v)$ is an E-mapping of $D^{r}$ satisfying

$$
f(0)=f_{0}(0) \quad \text { and } \quad f^{\prime}(0)=\lambda v \quad \text { for some } \lambda>0 .
$$

Proof. Observe that Proposition 5.1.1 provides us with a mapping $g_{0}=$ $\Phi \circ f_{0}$ and a domain $G_{0}:=\Phi\left(D_{0}\right)$ giving a data for situation ( $\dagger$ ) (here $\partial D_{0}$ is contained in $U_{0}$ ). Clearly, $\rho_{0}:=r_{0} \circ \Phi^{-1}$ is a defining function of $G_{0}$.

Using Lemmas 5.3.1 and 5.3.2, we get neighborhoods $V_{0}$ of $\left(\rho_{0}, g_{0}^{\prime}(0)\right)$ and $W_{0}$ of $\left(\rho_{0}, g_{0}\left(\xi_{0}\right)\right)$ and real-analytic mappings $\Upsilon, \widetilde{\Upsilon}$ such that $\Xi(\rho, v, \Upsilon(\rho, v))$ $=0$ on $V_{0}$ and $\widetilde{\Xi}(\rho, w, \widetilde{\Upsilon}(\rho, w))=0$ on $W_{0}$. Define

$$
\widehat{\Lambda}:=\pi_{B} \circ \widetilde{\Upsilon}, \quad \Omega:=\pi_{\mathbb{R}} \circ \widetilde{\Upsilon}, \quad \widehat{\Gamma}:=\pi_{B} \circ \Upsilon,
$$

where

$$
\pi_{B}: B \times Q_{0} \times \mathbb{R} \rightarrow B, \quad \pi_{\mathbb{R}}: B \times Q_{0} \times \mathbb{R} \rightarrow \mathbb{R}
$$

are the projections.
If $\rho$ is sufficiently close to $\rho_{0}$ then the hypersurface $\{\rho=0\}$ bounds a strongly linearly convex domain. Moreover, $\widehat{\Lambda}(\rho, w)$ and $\widehat{\Gamma}(\rho, v)$ are weak stationary mappings of $G^{\rho}$ (see Remark 5.5.1).

Composing $\widehat{\Lambda}(\rho, w)$ and $\widehat{\Gamma}(\rho, v)$ with $\Phi^{-1}$ and making use of Remark 5.1.2, we get weak stationary mappings of $D^{r}$, where $r:=\rho \circ \Phi$. To show that they are $E$-mappings we proceed as follows. If $D^{r}$ is sufficiently close to $D_{0}$ (this depends on the distance between $\rho$ and $\rho_{0}$ ) then the domain $D^{r}$ is strongly linearly convex, so by the results of Section 2 ,

$$
\Lambda(r, w):=\Phi^{-1} \circ \widehat{\Lambda}(\rho, w) \quad \text { and } \quad \Gamma(r, v):=\Phi^{-1} \circ \widehat{\Gamma}(\rho, v)
$$

are stationary mappings. Further, they are close to $f_{0}$ provided that $r$ is sufficiently close to $r_{0}$. Therefore, their winding numbers are equal. Thus $\Lambda(r, w)$ and $\Gamma(r, v)$ satisfy the condition (4) of Definition 1.12 , i.e. they are $E$-mappings.

## 6. Localization property

Proposition 6.1. Let $D \subset \mathbb{C}^{n}$, $n \geq 2$, be a domain such that $\partial D$ is real-analytic and strongly convex in a neighborhood of $a \in \partial D$. Then for any sufficiently small neighborhood $V_{0}$ of a there is a weak stationary mapping $f$ of $D \cap V_{0}$ such that $f(\mathbb{T}) \subset \partial D$. In particular, $f$ is a weak stationary mapping of $D$.

Proof. Let $r$ be a real-analytic defining function in a neighborhood of $a$. The problem we are dealing with has a local character, so replacing $r$ with $r \circ \Psi$, where $\Psi$ is a local biholomorphism near $a$, we may assume that $a=$ $(0, \ldots, 0,1)$ and a defining function of $D$ near $a$ is $r(z)=-1+|z|^{2}+h(z-a)$, where $h$ is real-analytic in a neighborhood of 0 and $h(z)=O\left(|z|^{3}\right)$ as $z \rightarrow 0$ (cf. [8, p. 321]).

Following [5], let us consider the mappings
$A_{t}(z):=\left(\left(1-t^{2}\right)^{1 / 2} \frac{z^{\prime}}{1+t z_{n}}, \frac{z_{n}+t}{1+t z_{n}}\right), \quad z=\left(z^{\prime}, z_{n}\right) \in \mathbb{C}^{n-1} \times \mathbb{D}, t \in(0,1)$, which restricted to $\mathbb{B}_{n}$ are automorphisms. Let

$$
r_{t}(z):= \begin{cases}\frac{\left|1+t z_{n}\right|^{2}}{1-t^{2}} r\left(A_{t}(z)\right), & t \in(0,1) \\ -1+|z|^{2}, & t=1\end{cases}
$$

It is clear that $f_{(1)}(\zeta)=(\zeta, 0, \ldots, 0), \zeta \in \mathbb{D}$, is a stationary mapping of $\mathbb{B}_{n}$. We want to have situation $(\dagger)$, which will allow us to use Lemma 5.3.1 (or Lemma 5.3.2). Note that $r_{t}$ does not converge to $r_{1}$ as $t \rightarrow 1$. However, $r_{t} \rightarrow r_{1}$ in $X_{0}\left(U_{0}, U_{0}^{\mathbb{C}}\right)$, where $U_{0}$ is a neighborhood of $f_{(1)}(\mathbb{T})$ contained in $\left\{z \in \mathbb{C}^{n}: \operatorname{Re} z_{n}>-1 / 2\right\}$ and $U_{0}^{\mathbb{C}}$ is sufficiently small (remember that $\left.h(z)=O\left(|z|^{3}\right)\right)$.

Therefore, making use of Lemma 5.3.1 for $t$ sufficiently close to 1 , we obtain stationary mappings $f_{(t)}$ of $D_{t}:=\left\{z \in \mathbb{C}^{n}: r_{t}(z)<0\right.$, $\left.\operatorname{Re} z_{n}>-1 / 2\right\}$ such that $f_{(t)} \rightarrow f_{(1)}$ in the $W^{2,2}$-norm (so also in the sup-norm). Actually, it follows from Lemma 5.3.1 that one may take $f_{(t)}:=\pi_{B} \circ \Upsilon\left(r_{t}, f_{(1)}^{\prime}(0)\right)$ (keeping the notation from that lemma). The argument used in Remark 5.5.1 shows that $f_{(t)}$ satisfies conditions $\left(1^{\prime}\right),\left(2^{\prime}\right)$ and $\left(3^{\prime}\right)$ of Definition 1.11. Since the non-constant function $r \circ A_{t} \circ f_{(t)}$ is subharmonic on $\mathbb{D}$, continuous on $\overline{\mathbb{D}}$, and $r \circ A_{t} \circ f_{(t)}=0$ on $\mathbb{T}$, we see from the maximum principle that $f_{(t)}$ maps $\mathbb{D}$ in $D_{t}$. Therefore, $f_{(t)}$ are weak stationary mappings for $t$ close to 1 .

In particular,

$$
f_{(t)}(\mathbb{D}) \subset 2 \mathbb{B}_{n} \cap\left\{z \in \mathbb{C}^{n}: \operatorname{Re} z_{n}>-1 / 2\right\}
$$

provided that $t$ is close to 1 . The mappings $A_{t}$ have the following important property:

$$
A_{t}\left(2 \mathbb{B}_{n} \cap\left\{z \in \mathbb{C}^{n}: \operatorname{Re} z_{n}>-1 / 2\right\}\right) \rightarrow\{a\}, \quad t \rightarrow 1
$$

in the sense of the Hausdorff distance.
Therefore, we find from Remark 5.1.2 that $g_{(t)}:=A_{t} \circ f_{(t)}$ is a stationary mapping of $D$. Since $g_{(t)}$ maps $\mathbb{D}$ onto an arbitrarily small neighborhood of $a$ provided that $t$ is sufficiently close to 1 , we immediately get the assertion. -
7. Proofs of Theorems 1.3 and 1.14. We start this section with the following

Lemma 7.1. For any different $z, w \in D$ (resp. for any $\left.z \in D, v \in\left(\mathbb{C}^{n}\right)_{*}\right)$ there exists an E-mapping $f: \mathbb{D} \rightarrow D$ such that $f(0)=z, f(\xi)=w$ for some $\xi \in(0,1)$ (resp. $f(0)=z, f^{\prime}(0)=\lambda v$ for some $\left.\lambda>0\right)$.

Proof. Fix different $z, w \in D$ (resp. $z \in D, v \in\left(\mathbb{C}^{n}\right)_{*}$ ).
First, consider the case when $D$ is bounded strongly convex with realanalytic boundary. Without loss of generality one may assume that $0 \in D \Subset$ $\mathbb{B}_{n}$. We need some properties of the Minkowski functionals.

Let $\mu_{G}$ be the Minkowski functional of a domain $G \subset \mathbb{C}^{n}$ containing the origin, i.e.

$$
\mu_{G}(x):=\inf \{s>0: x / s \in G\}, \quad x \in \mathbb{C}^{n}
$$

Assume that $G$ is bounded strongly convex with real-analytic boundary. We shall show that

- $\mu_{G}-1$ is a defining function of $G$, real-analytic outside 0 ;
- $\mu_{G}^{2}-1$ is a defining function of $G$, real-analytic outside 0 and strongly convex outside 0 .

Clearly, $G=\left\{x \in \mathbb{R}^{2 n}: \mu_{G}(x)<1\right\}$. Setting

$$
q(x, s):=r(x / s), \quad(x, s) \in U_{0} \times U_{1}
$$

where $r$ is a real-analytic defining function of $G$ and $U_{0} \subset \mathbb{R}^{2 n}, U_{1} \subset \mathbb{R}$ are neighborhoods of $\partial G$ and 1 respectively, we have

$$
\frac{\partial q}{\partial s}(x, s)=-\frac{1}{s^{2}}\left\langle\nabla r\left(\frac{x}{s}\right), x\right\rangle_{\mathbb{R}} \neq 0
$$

for $(x, s)$ such that $x \in \partial G$ and $s=\mu_{G}(x)=1$ (since $0 \in G$, the vector $-x$ hooked at the point $x$ points to the inside of $G$, so it is not orthogonal to the normal vector at $x$ ). By the Implicit Function Theorem for the equation
$q=0$, the function $\mu_{G}$ is real-analytic in a neighborhood $V_{0}$ of $\partial G$. To see that $\mu_{G}$ is also real-analytic outside 0 , fix $x_{0} \in\left(\mathbb{R}^{2 n}\right)_{*}$. Then the set

$$
W_{0}:=\left\{x \in \mathbb{R}^{2 n}: x / \mu_{G}\left(x_{0}\right) \in V_{0}\right\}
$$

is open and contains $x_{0}$. Since

$$
\mu_{G}(x)=\mu_{G}\left(x_{0}\right) \mu_{G}\left(\frac{x}{\mu_{G}\left(x_{0}\right)}\right), \quad x \in W_{0}
$$

the function $\mu_{G}$ is real-analytic in $W_{0}$. Therefore, we can apply $d / d s$ to both sides of $\mu_{G}(s x)=s \mu_{G}(x), x \neq 0, s>0$, to obtain

$$
\left\langle\nabla \mu_{G}(x), x\right\rangle_{\mathbb{R}}=\mu_{G}(x), \quad x \neq 0
$$

so $\nabla \mu_{G} \neq 0$ in $\left(\mathbb{R}^{2 n}\right)_{*}$.
Furthermore, $\nabla \mu_{G}^{2}=2 \mu_{G} \nabla \mu_{G}$, so $\mu_{G}^{2}-1$ is also a defining function of $G$. To show that $u:=\mu_{G}^{2}$ is strongly convex outside 0 , let us prove that

$$
X^{T} \mathcal{H}_{a} X>0, \quad a \in \partial G, X \in\left(\mathbb{R}^{2 n}\right)_{*},
$$

where $\mathcal{H}_{x}:=\mathcal{H} u(x)$ for $x \in\left(\mathbb{R}^{2 n}\right)_{*}$. Applying $\partial / \partial x_{j}$ to both sides of

$$
u(s x)=s^{2} u(x), \quad x, s \neq 0
$$

we get

$$
\begin{equation*}
\frac{\partial u}{\partial x_{j}}(s x)=s \frac{\partial u}{\partial x_{j}}(x) \tag{41}
\end{equation*}
$$

and further applying $d / d s$ we obtain

$$
\sum_{k=1}^{2 n} \frac{\partial^{2} u}{\partial x_{j} \partial x_{k}}(s x) x_{k}=\frac{\partial u}{\partial x_{j}}(x)
$$

In particular,

$$
x^{T} \mathcal{H}_{x} y=\sum_{j, k=1}^{2 n} \frac{\partial^{2} u}{\partial x_{k} \partial x_{j}}(x) x_{k} y_{j}=\langle\nabla u(x), y\rangle_{\mathbb{R}}, \quad x \in\left(\mathbb{R}^{2 n}\right)_{*}, y \in \mathbb{R}^{2 n}
$$

Let $a \in \partial G$. Since $\left\langle\nabla \mu_{G}(a), a\right\rangle_{\mathbb{R}}=\mu_{G}(a)=1$, we have $a \notin T_{G}^{\mathbb{R}}(a)$. Any $X \in\left(\mathbb{R}^{2 n}\right)_{*}$ can be represented as $\alpha a+\beta Y$, where $Y \in T_{G}^{\mathbb{R}}(a), \alpha, \beta \in \mathbb{R}$, $(\alpha, \beta) \neq(0,0)$. Then

$$
\begin{aligned}
X^{T} \mathcal{H}_{a} X & =\alpha^{2} a^{T} \mathcal{H}_{a} a+2 \alpha \beta a^{T} \mathcal{H}_{a} Y+\beta^{2} Y^{T} \mathcal{H}_{a} Y \\
& =\alpha^{2}\langle\nabla u(a), a\rangle_{\mathbb{R}}+2 \alpha \beta\langle\nabla u(a), Y\rangle_{\mathbb{R}}+\beta^{2} Y^{T} \mathcal{H}_{a} Y \\
& =\alpha^{2} 2 \mu_{G}(a)\left\langle\nabla \mu_{G}(a), a\right\rangle_{\mathbb{R}}+\beta^{2} Y^{T} \mathcal{H}_{a} Y=2 \alpha^{2}+\beta^{2} Y^{T} \mathcal{H}_{a} Y .
\end{aligned}
$$

Since $G$ is strongly convex, the Hessian of any defining function is strictly positive on the real tangent space, i.e. $Y^{T} \mathcal{H}_{a} Y>0$ if $Y \in\left(T_{G}^{\mathbb{R}}(a)\right)_{*}$. Hence $X^{T} \mathcal{H}_{a} X \geq 0$. Note that $X^{T} \mathcal{H}_{a} X=0$ cannot happen, since then $\alpha=0$,
consequently $\beta \neq 0$ and $Y^{T} \mathcal{H}_{a} Y=0$. On the other hand $Y=X / \beta \neq 0$, a contradiction.

Applying $\partial / \partial x_{k}$ to both sides of (41), we obtain

$$
\frac{\partial^{2} u}{\partial x_{j} \partial x_{k}}(s x)=\frac{\partial^{2} u}{\partial x_{j} \partial x_{k}}(x), \quad x, s \neq 0
$$

and for $a, X \in\left(\mathbb{R}^{2 n}\right)_{*}$,

$$
X^{T} \mathcal{H}_{a} X=X^{T} \mathcal{H}_{a / \mu_{G}(a)} X>0
$$

Let us consider the sets

$$
D_{t}:=\left\{x \in \mathbb{C}^{n}: t \mu_{D}^{2}(x)+(1-t) \mu_{\mathbb{B}_{n}}^{2}(x)<1\right\}, \quad t \in[0,1]
$$

Since the functions $t \mu_{D}^{2}+(1-t) \mu_{\mathbb{B}_{n}}^{2}$ are real-analytic in $\left(\mathbb{C}^{n}\right)_{*}$ and strongly convex in $\left(\mathbb{C}^{n}\right)_{*}, D_{t}$ are strongly convex domains with real-analytic boundaries satisfying

$$
D=D_{1} \Subset D_{t_{2}} \Subset D_{t_{1}} \Subset D_{0}=\mathbb{B}_{n} \quad \text { if } 0<t_{1}<t_{2}<1
$$

It is clear that $\mu_{D_{t}}=\sqrt{t \mu_{D}^{2}+(1-t) \mu_{\mathbb{B}_{n}}^{2}}$. Further, if $t_{1}$ is close to $t_{2}$ then $D_{t_{1}}$ is close to $D_{t_{2}}$ with respect to the topology introduced in Section 5. We want to show that $D_{t}$ are in some family $\mathcal{D}(c)$. Only the interior and exterior ball conditions need verification.

There exists $\delta>0$ such that $\delta \mathbb{B}_{n} \Subset D$. Further, $\nabla \mu_{D_{t}}^{2} \neq 0$ in $\left(\mathbb{R}^{2 n}\right)_{*}$. Set

$$
M:=\sup \left\{\frac{\mathcal{H} \mu_{D_{t}}^{2}(x ; X)}{\left|\nabla \mu_{D_{t}}^{2}(y)\right|}: t \in[0,1], x, y \in 2 \overline{\mathbb{B}}_{n} \backslash \delta \mathbb{B}_{n}, X \in \mathbb{R}^{2 n},|X|=1\right\}
$$

It is a positive number since the functions $\mu_{D_{t}}^{2}$ are strongly convex in $\left(\mathbb{R}^{2 n}\right)_{*}$ and the 'sup' of the continuous, positive function is taken over a compact set. Let

$$
r:=\min \left\{\frac{1}{2 M}, \frac{\operatorname{dist}\left(\partial D, \delta \mathbb{B}_{n}\right)}{2}\right\}
$$

For fixed $t \in[0,1]$ and $a \in \partial D_{t}$ put $a^{\prime}:=a-r \nu_{D_{t}}(a)$. In particular, $\overline{B_{n}\left(a^{\prime}, r\right)} \subset 2 \overline{\mathbb{B}}_{n} \backslash \delta \mathbb{B}_{n}$. Let us define

$$
h(x):=\mu_{D_{t}}^{2}(x)-\frac{\left|\nabla \mu_{D_{t}}^{2}(a)\right|}{2\left|a-a^{\prime}\right|}\left(\left|x-a^{\prime}\right|^{2}-r^{2}\right), \quad x \in 2 \overline{\mathbb{B}}_{n} \backslash \delta \mathbb{B}_{n}
$$

We have $h(a)=1$ and

$$
\nabla h(x)=\nabla \mu_{D_{t}}^{2}(x)-\frac{\left|\nabla \mu_{D_{t}}^{2}(a)\right|}{\left|a-a^{\prime}\right|}\left(x-a^{\prime}\right)
$$

For $x=a$, dividing the right side by $\left|\nabla \mu_{D_{t}}^{2}(a)\right|$, we get the difference of the same normal vectors $\nu_{D_{t}}(a)$, so $\nabla h(a)=0$. Moreover, for $|X|=1$,

$$
\mathcal{H} h(x ; X)=\mathcal{H} \mu_{D_{t}}^{2}(x ; X)-\frac{\left|\nabla \mu_{D_{t}}^{2}(a)\right|}{r} \leq M\left|\nabla \mu_{D_{t}}^{2}(a)\right|-2 M\left|\nabla \mu_{D_{t}}^{2}(a)\right|<0
$$

It follows that $h \leq 1$ in any convex set $S$ such that $a \in S \subset 2 \overline{\mathbb{B}}_{n} \backslash \delta \mathbb{B}_{n}$. Indeed, assume the contrary. Then there is $y \in S$ such that $h(y)>1$. Let us join $a$ and $y$ with a line segment and consider the function

$$
g:[0,1] \ni t \mapsto h(t a+(1-t) y) \in \mathbb{R}
$$

Since $a$ is a strong local maximum point of $h$, the function $g$ has a local minimum at some point $t_{0} \in(0,1)$. Hence

$$
0 \leq g^{\prime \prime}\left(t_{0}\right)=\mathcal{H} h\left(t_{0} a+\left(1-t_{0}\right) y ; a-y\right)
$$

which is impossible.
Setting $S:=\overline{B_{n}\left(a^{\prime}, r\right)}$, we get

$$
\mu_{D_{t}}^{2}(x) \leq 1+\frac{\left|\nabla \mu_{D_{t}}^{2}(a)\right|}{2\left|a-a^{\prime}\right|}\left(\left|x-a^{\prime}\right|^{2}-r^{2}\right)<1
$$

for $x \in B_{n}\left(a^{\prime}, r\right)$, i.e. $x \in D_{t}$.
The proof of the exterior ball condition is similar. Set

$$
m:=\inf \left\{\frac{\mathcal{H} \mu_{D_{t}}^{2}(x ; X)}{\left|\nabla \mu_{D_{t}}^{2}(y)\right|}: t \in[0,1], x, y \in\left(\overline{\mathbb{B}}_{n}\right)_{*}, X \in \mathbb{R}^{2 n},|X|=1\right\}
$$

Note that $m>0$. Actually, the homogeneity of $\mu_{D_{t}}$ implies $\mathcal{H} \mu_{D_{t}}^{2}(s x ; X)=$ $\mathcal{H} \mu_{D_{t}}^{2}(x ; X)$ and $\nabla \mu_{D_{t}}^{2}(s x)=s \nabla \mu_{D_{t}}^{2}(x)$ for $x \neq 0, X \in \mathbb{R}^{2 n}, s>0$. Therefore, there are positive constants $C_{1}, C_{2}$ such that $C_{1} \leq \mathcal{H} \mu_{D_{t}}^{2}(x ; X)$ for $x \neq 0, X \in \mathbb{R}^{2 n},|X|=1$ and $\left|\nabla \mu_{D_{t}}^{2}(y)\right| \leq C_{2}$ for $y \in \overline{\mathbb{B}}_{n}$. In particular, $m \geq C_{1} / C_{2}$.

Let $R:=2 / m$. For fixed $t \in[0,1]$ and $a \in \partial D_{t}$ put $a^{\prime \prime}:=a-R \nu_{D_{t}}(a)$. Let us define

$$
\widetilde{h}(x):=\mu_{D_{t}}^{2}(x)-\frac{\left|\nabla \mu_{D_{t}}^{2}(a)\right|}{2\left|a-a^{\prime \prime}\right|}\left(\left|x-a^{\prime \prime}\right|^{2}-R^{2}\right), \quad x \in \overline{\mathbb{B}}_{n}
$$

We have $\widetilde{h}(a)=1$ and

$$
\nabla \widetilde{h}(x)=\nabla \mu_{D_{t}}^{2}(x)-\frac{\left|\nabla \mu_{D_{t}}^{2}(a)\right|}{\left|a-a^{\prime \prime}\right|}\left(x-a^{\prime \prime}\right)
$$

so $\nabla \widetilde{h}(a)=0$. Moreover, for $x \in\left(\overline{\mathbb{B}}_{n}\right)_{*}$ and $|X|=1$,

$$
\mathcal{H} \widetilde{h}(x ; X)=\mathcal{H} \mu_{D_{t}}^{2}(x ; X)-\frac{\left|\nabla \mu_{D_{t}}^{2}(a)\right|}{R} \geq m\left|\nabla \mu_{D_{t}}^{2}(a)\right|-m\left|\nabla \mu_{D_{t}}^{2}(a)\right| / 2>0
$$

Therefore, $a$ is a strong local minimum point of $\widetilde{h}$.
Now using the properties listed above we may deduce that $\widetilde{h} \geq 1$ in $\overline{\mathbb{B}}_{n}$. We proceed as before: seeking a contradiction suppose that there is $y \in \overline{\mathbb{B}}_{n}$ such that $\widetilde{h}(y)<1$. Moving $y$ a little (if necessary) we may assume that 0 does not lie on the segment joining $a$ and $y$. Then the function $\widetilde{g}(t):=$ $\widetilde{h}(t a+(1-t) y)$ attains its local maximum at some point $t_{0} \in(0,1)$. The second derivative of $\widetilde{g}$ at $t_{0}$ is non-positive, which gives a contradiction with the positivity of the Hessian of $\widetilde{h}$.

Hence, we get

$$
\frac{\left|\nabla \mu_{D_{t}}^{2}(a)\right|}{2\left|a-a^{\prime \prime}\right|}\left(\left|x-a^{\prime \prime}\right|^{2}-R^{2}\right) \leq \mu_{D_{t}}^{2}(x)-1<0
$$

for $x \in D_{t}$, so $D_{t} \subset B_{n}\left(a^{\prime \prime}, R\right)$.
Let $T$ be the set of all $t \in[0,1]$ such that there is an $E$-mapping $f_{t}$ : $\mathbb{D} \rightarrow D_{t}$ with $f_{t}(0)=z, f_{t}\left(\xi_{t}\right)=w$ for some $\xi_{t} \in(0,1)$ (resp. $f_{t}(0)=z$, $f_{t}^{\prime}(0)=\lambda_{t} v$ for some $\left.\lambda_{t}>0\right)$. We claim that $T=[0,1]$. To prove this we will use the open-closed argument.

Clearly, $T \neq \emptyset$, as $0 \in T$. Moreover, $T$ is open in $[0,1]$. Indeed, let $t_{0} \in T$. It follows from Proposition 5.5.2 that there is a neighborhood $T_{0}$ of $t_{0}$ such that there exist $E$-mappings $f_{t}: \mathbb{D} \rightarrow D_{t}$ and $\xi_{t} \in(0,1)$ with $f_{t}(0)=z$, $f_{t}\left(\xi_{t}\right)=w$ for all $t \in T_{0}$ (resp. $\lambda_{t}>0$ such that $f_{t}(0)=z, f_{t}^{\prime}(0)=\lambda_{t} v$ for all $\left.t \in T_{0}\right)$.

To prove that $T$ is closed, choose a sequence $\left\{t_{m}\right\} \subset T$ converging to some $t \in[0,1]$. We want to show that $t \in T$. Since $f_{t_{m}}$ are $E$-mappings, they are complex geodesics. Therefore, making use of the inclusions $D \subset$ $D_{t_{m}} \subset \mathbb{B}_{n}$, we find that there is a compact set $K \subset(0,1)$ (resp. a compact set $\widetilde{K} \subset(0, \infty))$ such that $\left\{\xi_{t_{m}}\right\} \subset K$ (resp. $\left.\left\{\lambda_{t_{m}}\right\} \subset \widetilde{K}\right)$. By Propositions 4.7 and 4.9 the functions $f_{t_{m}}$ and ${\widetilde{f_{t}}}_{t_{m}}$ are equicontinuous in $\mathcal{C}^{1 / 2}(\overline{\mathbb{D}})$, and by Propositions 4.6 and 4.8 the functions $\rho_{t_{m}}$ are uniformly bounded from above and below by positive numbers and equicontinuous in $\mathcal{C}^{1 / 2}(\mathbb{T})$. From the Arzelà-Ascoli Theorem there are a subsequence $\left\{s_{m}\right\} \subset\left\{t_{m}\right\}$ and mappings $f, \widetilde{f} \in \mathcal{O}(\mathbb{D}) \cap \mathcal{C}^{1 / 2}(\overline{\mathbb{D}}), \rho \in \mathcal{C}^{1 / 2}(\mathbb{T})$ such that $f_{s_{m}} \rightarrow f, \widetilde{f}_{s_{m}} \rightarrow \widetilde{f}$ uniformly on $\overline{\mathbb{D}}, \rho_{s_{m}} \rightarrow \rho$ uniformly on $\mathbb{T}$ and $\xi_{s_{m}} \rightarrow \xi \in(0,1)$ (resp. $\lambda_{s_{m}} \rightarrow \lambda>0$ ).

Clearly, $f(\overline{\mathbb{D}}) \subset \bar{D}_{t}, f(\mathbb{T}) \subset \partial D_{t}$ and $\rho>0$. By the strong pseudoconvexity of $D_{t}$ we get $f(\mathbb{D}) \subset D_{t}$.

The conditions (3') and (4) of Definitions 1.11 and 1.12 follow from the uniform convergence of suitable functions. Therefore, $f$ is a weak $E$-mapping of $D_{t}$, and hence an $E$-mapping of $D_{t}$, satisfying $f(0)=z, f(\xi)=w$ (resp. $\left.f(0)=z, f^{\prime}(0)=\lambda v\right)$.

Let us go back to the general situation, that is, when a domain $D$ is bounded strongly linearly convex with real-analytic boundary. Take $\eta \in \partial D$ such that $\max _{\zeta \in \partial D}|z-\zeta|=|z-\eta|$. Then $\eta$ is a point of strong convexity of $D$. Indeed, by the Implicit Function Theorem one can assume that in a neighborhood of $\eta$ the defining functions of $D$ and $B:=B_{n}(z,|z-\eta|)$ are of the form $r(x):=\widetilde{r}(\widetilde{x})-x_{2 n}$ and $q(x):=\widetilde{q}(\widetilde{x})-x_{2 n}$ respectively, where $x=\left(\widetilde{x}, x_{2 n}\right) \in \mathbb{R}^{2 n}$ is sufficiently close to $\eta$. From the inclusion $D \subset B$ it follows that $r-q \geq 0$ near $\eta$ and $(r-q)(\eta)=0$. Thus the Hessian $\mathcal{H}(r-q)(\eta)$ is weakly positive on $\mathbb{C}^{n}$. Since $\mathcal{H} q(\eta)$ is strictly positive on $T_{B}^{\mathbb{R}}(\eta)_{*}=T_{D}^{\mathbb{R}}(\eta)_{*}$, we find that $\mathcal{H} r(\eta)$ is strictly positive on $T_{D}^{\mathbb{R}}(\eta)_{*}$ as well.

By a continuity argument, there is a convex neighborhood $V_{0}$ of $\eta$ such that all points from $\partial D \cap V_{0}$ are points of strong convexity of $D$. It follows from Proposition 6.1 (after shrinking $V_{0}$ if necessary) that there is a weak stationary mapping $g: \mathbb{D} \rightarrow D \cap V_{0}$ such that $g(\mathbb{T}) \subset \partial D$. In particular, $g$ is a weak stationary mapping of $D$. Since $D \cap V_{0}$ is convex, the condition on the winding number is satisfied on $D \cap V_{0}$ (and hence on the whole $D$ ). Consequently, $g$ is an $E$-mapping of $D$.

If $z=g(0), w=g(\xi)$ for some $\xi \in \mathbb{D}$ (resp. $\left.z=g(0), v=g^{\prime}(0)\right)$ then there is nothing to prove. In the other case let us take curves $\alpha:[0,1] \rightarrow D$, $\beta:[0,1] \rightarrow D$ joining $g(0)$ and $z$, and $g(1 / 2)$ and $w$ (resp. $\alpha:[0,1] \rightarrow D$, $\beta:[0,1] \rightarrow\left(\mathbb{C}^{n}\right)_{*}$ joining $g(0)$ and $z$, and $g^{\prime}(0)$ and $\left.v\right)$. We may assume that the images of $\alpha$ and $\beta$ are disjoint. Let $T$ be the set of all $t \in[0,1]$ such that there is an $E$-mapping $g_{t}: \mathbb{D} \rightarrow D$ such that $g_{t}(0)=\alpha(t), g_{t}\left(\xi_{t}\right)=\beta(t)$ for some $\xi_{t} \in(0,1)$ (resp. $g_{t}(0)=\alpha(t), g_{t}^{\prime}(0)=\lambda_{t} \beta(t)$ for some $\lambda_{t}>0$ ). Again $T \neq \emptyset$, since $0 \in T$. Using the results of Section 4 just as before (but for one domain), we see that $T$ is closed.

Since $\widetilde{k}_{D}$ is symmetric, it follows from Proposition 5.5.2 (1) that the set $T$ is open in $[0,1]$ (first we move along $\alpha$, then by symmetry we move along $\beta$ ). Therefore, $g_{1}$ is the $E$-mapping for $z, w$.

In the case of $\kappa_{D}$ we change a point and then we change a direction. To be more precise, consider the set $S$ of all $s \in[0,1]$ such that there is an $E$-mapping $h_{s}: \mathbb{D} \rightarrow D$ with $h_{s}(0)=\alpha(s)$. Then $0 \in S$, by Proposition 5.5.2(1) the set $S$ is open in $[0,1]$ and, by results of Section 4 again, it is closed. Hence $S=[0,1]$. Now we may join $h_{1}^{\prime}(0)$ and $v$ with a curve $\gamma:[0,1] \rightarrow \mathbb{C}^{n}$. Let us define $R$ as the set of all $r \in[0,1]$ such that there is an $E$-mapping $\widetilde{h}_{r}: \mathbb{D} \rightarrow D$ with $\widetilde{h}_{r}(0)=h_{1}(0), \widetilde{h}_{r}^{\prime}(0)=\sigma_{r} \gamma(1-r)$ for some $\sigma_{r}>0$. Then $1 \in R$, by Proposition 5.5.2(2) the set $R$ is open in [0,1], and by Section 4 it is closed. Hence $R=[0,1]$, so $\widetilde{h}_{0}$ is the $E$-mapping for $z, v$.

Now we are in a position to prove the main results of Lempert's paper.
Proof of Theorem 1.3 (real-analytic case). It follows from Lemma 7.1 that for any different points $z, w \in D$ (resp. $\left.z \in D, v \in\left(\mathbb{C}^{n}\right)_{*}\right)$ one may find an $E$-mapping passing through them (resp. $f(0)=z, f^{\prime}(0)=v$ ). On the other hand, it follows from Proposition 3.1.1 that $E$-mappings have left inverses, so they are complex geodesics.

Proof of Theorem 1.14 (real-analytic case). This is a direct consequence of Lemma 7.1 and Corollary 3.1.5

## $\mathcal{C}^{2}$-smooth case

Lemma 7.2. Let $D \subset \mathbb{C}^{n}$, $n \geq 2$, be a bounded strongly pseudoconvex domain with $\mathcal{C}^{2}$-smooth boundary. Take $z \in D$ and let $r$ be a defining function of $D$ such that

- $r \in \mathcal{C}^{2}\left(\mathbb{C}^{n}\right)$;
- $D=\left\{x \in \mathbb{C}^{n}: r(x)<0\right\}$;
- $\mathbb{C}^{n} \backslash \bar{D}=\left\{x \in \mathbb{C}^{n}: r(x)>0\right\}$;
- $|\nabla r|=1$ on $\partial D$;
- $\sum_{j, k=1}^{n} \frac{\partial^{2} r}{\partial z_{j} \partial \bar{z}_{k}}(a) X_{j} \bar{X}_{k} \geq C|X|^{2}$ for any $a \in \partial D$ and $X \in \mathbb{C}^{n}$ with some constant $C>0$.

Suppose that there is a sequence $\left\{r_{m}\right\}$ of $\mathcal{C}^{2}$-smooth real-valued functions such that $\partial^{\alpha} r_{m} / \partial x^{\alpha} \rightarrow \partial^{\alpha} r / \partial x^{\alpha}$ locally uniformly for any $\alpha \in \mathbb{N}_{0}^{2 n}$ with $|\alpha| \leq 2$. Let $D_{m}$ be the connected component of the set $\left\{x \in \mathbb{C}^{n}: r_{m}(x)<0\right\}$ such that $D_{m}$ contains the point $z$.

Then there is $c>0$ such that $\left(D_{m}, z\right)$ and $(D, z)$ belong to $\mathcal{D}(c), m \gg 1$.
Proof. Losing no generality assume that $D \Subset \mathbb{B}_{n}$. Note that the conditions (1), (5), (6) of Definition 4.1 are clearly satisfied. To find $c$ satisfying (2), we take $s>0$ such that $\mathcal{H} r(x ; X)<s|X|^{2}$ for $x \in \overline{\mathbb{B}}_{n}$ and $X \in\left(\mathbb{R}^{2 n}\right)_{*}$. Then $\mathcal{H} r_{m}(x ; X)<2 s|X|^{2}$ for $x \in \overline{\mathbb{B}}_{n}, X \in\left(\mathbb{R}^{2 n}\right)_{*}$ and $m \gg 1$. Let $U_{0} \subset \mathbb{B}_{n}$ be an open neighborhood of $\partial D$ such that $|\nabla r|$ is on $U_{0}$ between $3 / 4$ and $5 / 4$. Note that $\partial D_{m} \subset U_{0}$ and $\left|\nabla r_{m}\right| \in(1 / 2,3 / 2)$ on $U_{0}$ for $m \gg 1$.

Fix $m$ and $a \in \partial D_{m}$ and put $b:=a-R \nu_{D_{m}}(a)$, where the small number $R>0$ will be specified later. There is $t>0$ such that $\nabla r_{m}(a)=2 t(a-b)$. Note that $t$ may be arbitrarily large provided that $R$ was small enough. We take $t:=2 s$ and $R:=\left|\nabla r_{m}(a)\right| / t$. Then we have $\mathcal{H} r_{m}(x ; X)<2 t|X|^{2}$ for $x \in \overline{\mathbb{B}}_{n}, X \in\left(\mathbb{R}^{2 n}\right)_{*}$ and $m \gg 1$. Hence the function

$$
h(x):=r_{m}(x)-t\left(|x-b|^{2}-R^{2}\right), \quad x \in \mathbb{C}^{n}
$$

attains at $a$ its global maximum on $\overline{\mathbb{B}}_{n}$ (since $a$ is a strong local maximum point and the Hessian of $h$ is negative on the convex set $\overline{\mathbb{B}}_{n}$, cf. the proof of Lemma 7.1). Thus $h \leq 0$ on $\mathbb{B}_{n}$. From this we immediately get (2).

Note that it follows from (2) that $D_{m}=\left\{x \in \mathbb{C}^{n}: r_{m}(x)<0\right\}$ for $m$ big enough (i.e. the set $\left\{x \in \mathbb{C}^{n}: r_{m}(x)<0\right\}$ is connected).

Moreover, the condition (2) implies the condition (3) as follows. We infer from Remark 4.4 that there is $c^{\prime}>0$ such that $D$ satisfies (3) with $c^{\prime}$. Let $m_{0}$ be such that the Hausdorff distance between $\partial D$ and $\partial D_{m}$ is smaller than $1 / c^{\prime}$ for $m \geq m_{0}$. There is $c^{\prime \prime}>0$ such that $D_{m_{0}}$ satisfies (3) with $c^{\prime \prime}$. Losing no generality we may assume that $c^{\prime \prime}<c^{\prime}$. Take any $x, y \in D_{m}$. Since $D_{m}$ satisfies (without loss of generality) the interior ball condition with radius $1 / c^{\prime}$, we infer that there are balls of radius $1 / c^{\prime}$ contained in $D_{m}$ and containing $x$ and $y$ respectively. The centers of these balls lie in $D_{m_{0}}$. Using the fact that $\left(D_{m_{0}}, z\right)$ lies in $\mathcal{D}\left(c^{\prime \prime}\right)$, we may join the chosen centers to balls of radius $1 /\left(2 c^{\prime \prime}\right)$ as in the condition (3), so we have found a chain consisting of balls of radii $1 / c^{\prime}$ and $1 / c^{\prime \prime}$, joining $x$ and $y$.

Thus we may join $x$ and $y$ with balls contained entirely in the constructed chain, of the same radii depending only on $c^{\prime}$ and $c^{\prime \prime}$.

Now we prove (4). We shall show that there is $c>c^{\prime}$ such that every $D_{m}$ satisfies (4) with $c$ for $m$ large enough. To do this let us cover $\partial D$ with a finite number of balls $B_{j}, j=1, \ldots, N$, from the condition (4) and let $B_{j}^{\prime}$ be a ball relatively contained in $B_{j}$ such that $\left\{B_{j}^{\prime}\right\}_{j=1}^{N}$ covers $\partial D$ as well. Let $\Phi_{j}$ be mappings corresponding to $B_{j}$. Let $\varepsilon$ be such that any ball of radius $\varepsilon$ intersecting $\partial D$ is relatively contained in $B_{j}^{\prime}$ for some $j$. Observe that any ball $B$ of radius $\varepsilon / 2$ intersecting $\partial D_{m}$ is contained in a ball of radius $\varepsilon$ intersecting $\partial D$; hence it is contained in $B_{j}^{\prime}$ for some $j$. Then the pair $B, \Phi_{j}$ satisfies the conditions $(4)(b, c, d)$. Therefore, it suffices to check that there is $c>2 / \varepsilon$ such that each pair $B_{j}^{\prime}, \Phi_{j}$ satisfies the condition (4) for $D_{m}$ with $c$ (and $m \gg 1$ ). This is possible since $\Phi_{j}\left(D_{m}\right) \subset \Phi_{j}(D), D^{\alpha} \Phi_{j}\left(\partial D_{m} \cap B_{j}\right)$ converges to $D^{\alpha} \Phi_{j}\left(\partial D \cap B_{j}\right)$ for $|\alpha| \leq 2$ and for any $w \in \Phi\left(\partial D \cap B_{j}\right)$ there is a ball of radius $2 / \varepsilon$ containing $\Phi_{j}(D)$ and tangent to $\partial \Phi_{j}(D)$ at $w$. To be precise, we proceed as follows.

Let $a, b \in \mathbb{C}^{n}$ and $x \in \partial B_{n}(a, \widetilde{c})$, where $\widetilde{c}>c^{\prime}$. Then a ball $B_{n}(2 a-x, 2 \widetilde{c})$ contains $B_{n}(a, \widetilde{c})$ and is tangent to $B_{n}(a, \widetilde{c})$ at $x$. There is a number $\eta=$ $\eta(\delta, \widetilde{c})>0$, independent of $a, b, x$, such that the diameter of the set $B_{n}(b, \widetilde{c}) \backslash$ $B_{n}(2 a-x, 2 \widetilde{c})$ is smaller than $\delta>0$ whenever $|a-b|<\eta$ (this is a simple consequence of the triangle inequality).

Let $\widetilde{s}>0$ be such that $\mathcal{H}\left(r \circ \Phi_{j}^{-1}\right)(x ; X) \geq 2 \widetilde{s}|X|^{2}$ for $x \in U_{j}, j=$ $1, \ldots, N$, where $U_{j}$ is an open neighborhood of $\Phi_{j}\left(\partial D \cap B_{j}\right)$. Then, for $m$ large enough, $\mathcal{H}\left(r_{m} \circ \Phi_{j}^{-1}\right)(x ; X) \geq \widetilde{s}|X|^{2}$ for $x \in U_{j}$ and $\Phi_{j}\left(\partial D_{m} \cap B_{j}^{\prime}\right) \subset U_{j}$, $j=1, \ldots, N$. Repeating for the function

$$
x \mapsto\left(r_{m} \circ \Phi_{j}^{-1}\right)(x)-\widetilde{t}\left(|x-\widetilde{b}|^{2}-\widetilde{R}^{2}\right)
$$

the argument used in the interior ball condition with suitably chosen $\widetilde{t}$ and uniform $\widetilde{R}>c$, we find that there is a uniform $\widetilde{\varepsilon}>0$ such that for any $j, m$ and $w \in \Phi_{j}\left(\partial D_{m} \cap B_{j}^{\prime}\right)$ there is a ball $B$ of radius $\widetilde{R}$, tangent to $\Phi_{j}\left(\partial D_{m} \cap B_{j}^{\prime}\right)$ at $w$, such that $\Phi_{j}\left(\partial D_{m} \cap B_{j}^{\prime}\right) \cap B_{n}(w, \widetilde{\varepsilon}) \subset B$. Let $a_{j, m}(w)$ denote its center.

On the other hand, for any $w \in \Phi_{j}\left(\partial D_{m} \cap B_{j}^{\prime}\right)$ there is $t>0$ such that $w^{\prime}=w+t \nu(w) \in \Phi_{j}\left(\partial D \cap B_{j}\right)$, where $\nu(w)$ is a normal vector to $\Phi_{j}\left(\partial D_{m} \cap B_{j}^{\prime}\right)$ at $w$. Let $a_{j}\left(w^{\prime}\right)$ be the center of a ball of radius $\widetilde{R}$ tangent to $\Phi_{j}\left(\partial D \cap B_{j}\right)$ at $w^{\prime}$. It follows that $\left|a_{j, m}(w)-a_{j}\left(w^{\prime}\right)\right|<\eta(\widetilde{\varepsilon} / 2, \widetilde{R})$ provided that $m$ is large enough.

Combining the facts presented above, we finish the proof of the exterior ball condition (with a radius dependent only on $\widetilde{\varepsilon}$ and $\widetilde{R}$ ).

Proof of Theorems 1.3 and 1.14 ( $\mathcal{C}^{2}$-smooth case). Losing no generality assume that $0 \in D \Subset \mathbb{B}_{n}$.

It follows from the Weierstrass Theorem that there is sequence $\left\{P_{k}\right\}$ of real polynomials on $\mathbb{C}^{n} \simeq \mathbb{R}^{2 n}$ such that

$$
\partial^{\alpha} P_{k} / \partial x^{\alpha} \rightarrow \partial^{\alpha} r / \partial x^{\alpha} \quad \text { uniformly on } \overline{\mathbb{B}}_{n}
$$

where $\alpha=\left(\alpha_{1}, \ldots, \alpha_{2 n}\right) \in \mathbb{N}_{0}^{2 n}$ is such that $|\alpha| \leq 2$. Consider the open set

$$
\widetilde{D}_{k, \varepsilon}:=\left\{x \in \mathbb{C}^{n}: P_{k}(x)+\varepsilon<0\right\} .
$$

Let $\varepsilon_{m}>0$ be a sequence converging to 0 with $3 \varepsilon_{m+1}<\varepsilon_{m}$. For any $m \in \mathbb{N}$ there is $k_{m} \in \mathbb{N}$ such that $\sup _{\overline{\mathbb{B}}_{n}}\left|P_{k_{m}}-r\right|<\varepsilon_{m}$. Putting $r_{m}:=P_{k_{m}}+2 \varepsilon_{m}$, we get $r+\varepsilon_{m}<r_{m}<r+3 \varepsilon_{m}$ on $\overline{\mathbb{B}}_{n}$. In particular, $r_{m+1}<r_{m}$ on $\overline{\mathbb{B}}_{n}$.

Let $D_{m}$ be the connected component of $D_{k_{m}, 2 \varepsilon_{m}}$ containing 0 . It is a bounded strongly linearly convex domain with real-analytic boundary and $r_{m}$ is its defining function provided that $m$ is large enough. Moreover, $D_{m} \subset D_{m+1}$ and $\bigcup_{m} D_{m}=D$. Using properties of holomorphically invariant functions and metrics we get Theorem 1.3 .

We are left with showing the claim that for any different $z, w \in D$ (resp. $z \in D, v \in\left(\mathbb{C}^{n}\right)_{*}$ ) there is a weak $E$-mapping for $z, w$ (resp. for $z, v$ ). Fix $z \in D$ and $w \in D\left(\right.$ resp. $\left.v \in\left(\mathbb{C}^{n}\right)_{*}\right)$. Then $z, w \in D_{m}\left(\right.$ resp. $\left.z \in D_{m}\right)$, $m \gg 1$. Therefore, for any $m \gg 1$ one may find an $E$-mapping $f_{m}$ of $D_{m}$ for $z, w($ resp. for $z, v)$. Since $\left(D_{m}, z\right) \in \mathcal{D}(c)$ for some uniform $c>0(m \gg 1)$ (Lemma 7.2), we find that $f_{m}, \widetilde{f}_{m}$ and $\rho_{m}$ satisfy the uniform estimates from Section 4. Thus, passing to a subsequence we may assume that $\left\{f_{m}\right\}$ converges uniformly on $\overline{\mathbb{D}}$ to a mapping $f \in \mathcal{O}(\mathbb{D}) \cap \mathcal{C}^{1 / 2}(\overline{\mathbb{D}})$ passing through $z, w$ (resp. such that $f(0)=z, f^{\prime}(0)=\lambda v, \lambda>0$ ), that $\left\{\widetilde{f}_{m}\right\}$ converges uniformly on $\overline{\mathbb{D}}$ to a mapping $\widetilde{f} \in \mathcal{O}(\mathbb{D}) \cap \mathcal{C}^{1 / 2}(\overline{\mathbb{D}})$ and that $\left\{\rho_{m}\right\}$ is convergent uniformly on $\mathbb{T}$ to a positive function $\rho \in \mathcal{C}^{1 / 2}(\mathbb{T})$ (in particular, $f^{\prime} \bullet \widetilde{f}=1$ in $\mathbb{D}$, so $\widetilde{f}$ has no zeroes in $\overline{\mathbb{D}})$. We already know that this implies that $f$ is a weak $E$-mapping of $D$.

To get $\mathcal{C}^{k-1-\varepsilon}$-smoothness of $f$ and its associated mappings for $k=$ $3,4, \ldots, \infty$, it suffices to repeat the proof of Proposition 5 in [5]. This is just the Webster Lemma (we have proved it in the real-analytic case-see Proposition 2.1). Namely, let

$$
\psi: \partial D \ni z \mapsto\left(z, T_{D}^{\mathbb{C}}(z)\right) \in \mathbb{C}^{n} \times\left(\mathbb{P}^{n-1}\right)_{*}
$$

where $\mathbb{P}^{n-1}$ is the $(n-1)$-dimensional complex projective space. Let $\pi$ : $\left(\mathbb{C}^{n}\right)_{*} \rightarrow \mathbb{P}^{n-1}$ be the canonical projection.

By the Webster Lemma, $\psi(\partial D)$ is a totally real manifold of class $\mathcal{C}^{k-1}$. Observe that the mapping $(f, \pi \circ \widetilde{f}): \overline{\mathbb{D}} \rightarrow \mathbb{C}^{n} \times \mathbb{P}^{n-1}$ is $1 / 2$-Hölder continuous, is holomorphic in $\mathbb{D}$ and maps $\mathbb{T}$ into $\psi(\partial D)$. Therefore, it is $\mathcal{C}^{k-1-\varepsilon_{-}}$ smooth for any $\varepsilon>0$, whence $f$ is $\mathcal{C}^{k-1-\varepsilon}$-smooth. Since $\nu_{D} \circ f$ is of class $\mathcal{C}^{k-1-\varepsilon}$, it suffices to proceed as in the proof of Proposition 2.1.

## 8. Appendix

### 8.1. Totally real submanifolds

Lemma 8.1.1. Let $M \subset \mathbb{C}^{m}$ be a totally real local $\mathcal{C}^{\omega}$ submanifold of real dimension $m$ and let $z \in M$. Then there are neighborhoods $U, V \subset \mathbb{C}^{m}$ of 0 and $z$ respectively and a biholomorphic mapping $\Phi: U \rightarrow V$ such that $\Phi\left(\mathbb{R}^{m} \cap U\right)=M \cap V$.

Lemma 8.1.2 (Reflection Principle). Let $M \subset \mathbb{C}^{m}$ be a totally real local $\mathcal{C}^{\omega}$ submanifold of real dimension $m$. Let $V_{0} \subset \mathbb{C}$ be a neighborhood of a point $\zeta_{0} \in \mathbb{T}$ and let $g: \overline{\mathbb{D}} \cap V_{0} \rightarrow \mathbb{C}^{m}$ be a continuous mapping. Suppose that $g \in \mathcal{O}\left(\mathbb{D} \cap V_{0}\right)$ and $g\left(\mathbb{T} \cap V_{0}\right) \subset M$. Then $g$ extends holomorphically past $\mathbb{T} \cap V_{0}$.
8.2. Some estimates of holomorphic functions of class $\mathcal{C}^{\alpha}$. Let us recall some theorems about functions holomorphic in $\mathbb{D}$ and continuous in $\overline{\mathbb{D}}$. Concrete values of the constants $M, K$ can be calculated by inspecting the proofs. In fact, it is only important that they do not depend on the functions involved.

Theorem 8.2.1 (Hardy, Littlewood, cf. [3, Theorem 3, p. 411]). Let $f \in$ $\mathcal{O}(\mathbb{D}) \cap \mathcal{C}(\overline{\mathbb{D}})$. Then for $\alpha \in(0,1]$ the following conditions are equivalent:

$$
\begin{align*}
& \exists M>0:\left|f\left(e^{i \theta}\right)-f\left(e^{i \theta^{\prime}}\right)\right| \leq M\left|\theta-\theta^{\prime}\right|^{\alpha}, \theta, \theta^{\prime} \in \mathbb{R} ;  \tag{42}\\
& \exists K>0:\left|f^{\prime}(\zeta)\right| \leq K(1-|\zeta|)^{\alpha-1}, \zeta \in \mathbb{D} \tag{43}
\end{align*}
$$

Moreover, if $M$ satisfying (42) is given then $K$ can be chosen as

$$
2^{(1-3 \alpha) / 2} \pi^{\alpha} M \int_{0}^{\infty} \frac{t^{\alpha}}{1+t^{2}} d t
$$

and if $K$ satisfying (43) is given then $M$ can be chosen as $(2 / \alpha+1) K$.
Theorem 8.2.2 (Hardy, Littlewood, cf. [3, Theorem 4, p. 413]). Let $f \in$ $\mathcal{O}(\mathbb{D}) \cap \mathcal{C}(\overline{\mathbb{D}})$ be such that

$$
\left|f\left(e^{i \theta}\right)-f\left(e^{i \theta^{\prime}}\right)\right| \leq M\left|\theta-\theta^{\prime}\right|^{\alpha}, \quad \theta, \theta^{\prime} \in \mathbb{R}
$$

for some $\alpha \in(0,1]$ and $M>0$. Then

$$
\left|f(\zeta)-f\left(\zeta^{\prime}\right)\right| \leq K\left|\zeta-\zeta^{\prime}\right|^{\alpha}, \quad \zeta, \zeta^{\prime} \in \overline{\mathbb{D}},
$$

where

$$
K:=\max \left\{2^{1-2 \alpha} \pi^{\alpha} M, 2^{(3-5 \alpha) / 2} \pi^{\alpha} \alpha^{-1} M \int_{0}^{\infty} \frac{t^{\alpha}}{1+t^{2}} d t\right\}
$$

Theorem 8.2.3 (Privalov, cf. [3, Theorem 5, p. 414]). Let $f \in \mathcal{O}(\mathbb{D})$ be such that $\operatorname{Re} f$ extends continuously on $\overline{\mathbb{D}}$ and

$$
\left|\operatorname{Re} f\left(e^{i \theta}\right)-\operatorname{Re} f\left(e^{i \theta^{\prime}}\right)\right| \leq M\left|\theta-\theta^{\prime}\right|^{\alpha}, \quad \theta, \theta^{\prime} \in \mathbb{R}
$$

for some $\alpha \in(0,1)$ and $M>0$. Then $f$ extends continuously on $\overline{\mathbb{D}}$ and

$$
\left|f(\zeta)-f\left(\zeta^{\prime}\right)\right| \leq K\left|\zeta-\zeta^{\prime}\right|^{\alpha}, \quad \zeta, \zeta^{\prime} \in \overline{\mathbb{D}},
$$

where

$$
\begin{aligned}
K:= & \max \left\{2^{1-2 \alpha} \pi^{\alpha}, 2^{(3-5 \alpha) / 2} \pi^{\alpha} \alpha^{-1} \int_{0}^{\infty} \frac{t^{\alpha}}{1+t^{2}} d t\right\} \\
& \times\left(\frac{2}{\alpha}+1\right) 2^{(3-3 \alpha) / 2} \pi^{\alpha} M \int_{0}^{\infty} \frac{t^{\alpha}}{1+t^{2}} d t
\end{aligned}
$$

8.3. Sobolev space. The Sobolev space $W^{2,2}(\mathbb{T})=W^{2,2}\left(\mathbb{T}, \mathbb{C}^{m}\right)$ is the space of functions $f: \mathbb{T} \rightarrow \mathbb{C}^{m}$, whose first two derivatives (in the sense of distributions) are in $L^{2}(\mathbb{T})$ (here we use the standard identification of functions on the unit circle and functions on $[0,2 \pi])$. Each such $f$ is $\mathcal{C}^{1}$ smooth.
$W^{2,2}(\mathbb{T})$ is a complex Hilbert space with the scalar product

$$
\langle f, g\rangle_{W}:=\langle f, g\rangle_{L}+\left\langle f^{\prime}, g^{\prime}\right\rangle_{L}+\left\langle f^{\prime \prime}, g^{\prime \prime}\right\rangle_{L},
$$

where

$$
\langle\widetilde{f}, \widetilde{g}\rangle_{L}:=\frac{1}{2 \pi} \int_{0}^{2 \pi}\left\langle\widetilde{f}\left(e^{i t}\right), \widetilde{g}\left(e^{i t}\right)\right\rangle d t .
$$

Let $\|\cdot\|_{L},\|\cdot\|_{W}$ denote the norms induced by $\langle\cdot,-\rangle_{L}$ and $\langle\cdot,-\rangle_{W}$. The following characterization follows easily from Parseval's identity:

$$
W^{2,2}(\mathbb{T})=\left\{f \in L^{2}(\mathbb{T}): \sum_{k=-\infty}^{\infty}\left(1+k^{2}+k^{4}\right)\left|a_{k}\right|^{2}<\infty\right\},
$$

where $a_{k} \in \mathbb{C}^{m}$ are the $m$-dimensional Fourier coefficients of $f$, i.e. $f(\zeta)=$ $\sum_{k=-\infty}^{\infty} a_{k} \zeta^{k}, \zeta \in \mathbb{T}$. More precisely, Parseval's identity gives

$$
\|f\|_{W}=\sqrt{\sum_{k=-\infty}^{\infty}\left(1+k^{2}+k^{4}\right)\left|a_{k}\right|^{2}}, \quad f \in W^{2,2}(\mathbb{T})
$$

Note that $W^{2,2}(\mathbb{T}) \subset \mathcal{C}^{1 / 2}(\mathbb{T}) \subset \mathcal{C}(\mathbb{T})$ and both inclusions are continuous (in particular, both inclusions are real-analytic). Note also that

$$
\|f\|_{\sup } \leq \sum_{k=-\infty}^{\infty}\left|a_{k}\right| \leq \sqrt{\sum_{k=-\infty}^{\infty} \frac{1}{1+k^{2}} \sum_{k=-\infty}^{\infty}\left(1+k^{2}\right)\left|a_{k}\right|^{2}} \leq \frac{\pi}{\sqrt{3}}\|f\|_{W}
$$

Moreover, direct calculation shows that there exists $C>0$ such that

$$
\left\|h^{\alpha}\right\|_{W} \leq C^{|\alpha|}\left\|h_{1}\right\|_{W}^{\alpha_{1}} \ldots \cdot\left\|h_{2 n}\right\|_{W}^{\alpha_{2 n}}, \quad h \in W^{2,2}\left(\mathbb{T}, \mathbb{C}^{n}\right), \alpha \in \mathbb{N}_{0}^{2 n}
$$

### 8.4. Matrices

Proposition 8.4.1 (Lempert, [5, Théorème B]). Let $A: \mathbb{T} \rightarrow \mathbb{C}^{n \times n}$ be a matrix-valued real-analytic mapping such that $A(\zeta)$ is self-adjoint and strictly positive for any $\zeta \in \mathbb{T}$. Then there exists $H \in \mathcal{O}\left(\overline{\mathbb{D}}, \mathbb{C}^{(n-1) \times(n-1)}\right)$ such that $\operatorname{det} H \neq 0$ on $\overline{\mathbb{D}}$ and $H H^{*}=A$ on $\mathbb{T}$.

Proposition 8.4.2 ([9, Lemma 2.1]). Let $A$ be a complex symmetric $n \times n$ matrix. Then

$$
\|A\|=\sup \left\{\left|z^{T} A z\right|: z \in \mathbb{C}^{n},|z|=1\right\}
$$

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