# QUASI-CONFORMAL MAPPINGS IN $n$-SPACE AND THE RIGIDITY OF HYPERBOLIC SPACE FORMS 

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## INTRODUCTION

The phenomenon that we wish to present in this paper can perhaps best be introduced by directing attention to the familiar case of compact Riemann surfaces.

Let $Y$ and $Y^{\prime}$ be diffeomorphic compact Riemann surfaces of genus greater than I . From the classical theory of uniformizing parameters, we may regard $Y$ and $Y^{\prime}$ as two dimensional Riemannian manifolds having constant negative curvature. We know moreover that Y and $\mathrm{Y}^{\prime}$ need not be conformally equivalent. A special case of our main theorem asserts (cf. Corollary 12.2):

Let Y and $\mathrm{Y}^{\prime}$ be diffeomorphic compact Riemannian manifolds having constant negative curvature. Then Y and $\mathrm{Y}^{\prime}$ are conformally equivalent, provided their dimension is greater than two.

The conformal mapping $\Psi$ of $Y$ to $Y^{\prime}$ can be chosen so that it is homotopic to any given diffeomorphism $\Phi$ of $Y$ to $Y^{\prime}$, and this conformal mapping $\Psi$ is unique.

Corollary 12.2 deals with the more general case that Y and $\mathrm{Y}^{\prime}$ have finite measure. To conclude that they are conformally equivalent, we have to make the hypothesis that there is a homeomorphism $\Phi: \mathrm{Y} \rightarrow \mathrm{Y}^{\prime}$ which is quasi-conformal. Indeed, the backbone of the method in this paper is the application of the theory of quasi-conformal mappings to the study of the following question:

Suppose $\Phi$ is a homeomorphism of $n$-dimensional hyperbolic space X onto itself, and suppose $\Gamma$ and $\Gamma^{\prime}$ are groups of isometries of X onto X . ( $\Gamma$ and $\Gamma^{\prime}$ correspond to the fundamental groups of Y and $\mathrm{Y}^{\prime}$.) Let $\theta: \Gamma \rightarrow \Gamma^{\prime}$ be an isomorphism and assume

$$
\begin{equation*}
\Phi(\gamma x)=\theta(\gamma) \Phi(x) \tag{*}
\end{equation*}
$$

for all $x \in \mathrm{X}, \gamma \in \Gamma$. Then $\Gamma^{\prime}=\Phi \Gamma \Phi^{-1}$ so that by hypothesis $\Gamma$ and $\Gamma^{\prime}$ are conjugate in the group of homeomorphisms of X onto X . Are $\Gamma$ and $\Gamma^{\prime}$ conjugate in the group $G$ of isometries of X ?

The central idea of our method is to study the effect of $\Phi$ " at infinity ". That is, if we look at the two dimensional case from the point of view of transformation groups, what distinguishes the case that $\Gamma$ and $\Gamma^{\prime}$ are conjugate from the contrary case? If we think of $X$ as the interior of the unit disc, and the compact Riemann surfaces as $Y=\Gamma \backslash X$ and $\mathrm{Y}^{\prime}=\Gamma^{\prime} \backslash \mathrm{X}$, then the map $\Phi$ will be smooth on X if $\Phi$ is a diffeomorphism, but it will not be smoothly extendible to the boundary unless $\Gamma$ is conjugate to $\Gamma^{\prime}$ in the group G of isometries.

This last assertion is valid more generally for arbitrary symmetric Riemannian spaces X (cf. Mostow, On the conjugacy of subgroups of semi-simple groups, Proc. of Symposia on Pure Math, v. 9 (1966), pp. 413-419) and that general fact motivates the viewpoint adopted here.

Indeed, onde our attention is directed to the behaviour of $\Phi$ at the boundary of X , the relevance of quasi-conformal mappings comes into view. The young Japanese mathematician, A. Mori, in a posthumous article: On quasi-conformality and pseudo-analyticity, Trans. Amer. Math. Soc., v. 84 (1957), pp. 56-77, proved in dimension 2 :

Theorem (10.1). - A quasi-conformal mapping of the open $n$-ball induces a homeomorphism on the boundary.

Mori's paper contains the seeds of not only the generalization to $n$ dimensions, but even of the more striking Theorem (ro.2) which states in effect:

A quasi-conformal mapping of the open $n$-ball induces a quasi-conformal mapping on the boundary. This result was first proved in dimension 3 by F. W. Gehring in "Rings and quasi-conformal mappings in space ", Trans. Amer. Math. Soc., v. io3 (1962), pp. 353-393. Unable, as was Mori, to use an inequality of Teichmüller on the moduli of rings because the proof rested on planar conformal mapping theory, Gehring achieved the requisite inequality in dimension 3 by proceeding from Loewner's " conformal capacity " (cf. Gehring, Symmetrization of rings in space, Trans. Amer. Math. Soc., v. ioi (1961), pp. 449-519). In this paper we develop the theory of quasi-conformal mappings in $n$-space adapting in large measure the procedure of Gehring.

The proof of the main theorem of this paper, Theorem (12.1), comes in the last section. It depends on an analysis of the boundary mapping induced by a quasiconformal mapping of hyperbolic $n$-space which satisfies condition (*) above. A key role in this analysis is played by an ergodic theorem for semi-groups in a semi-simple Lie group due essentially to F. I. Mautner. We have tried to reduce to the minimum the mathematics that our proof depends on, and thereby we have arrived at a relatively self-contained presentation. The expert on quasi-conformal mappings may recognize here and there some new proofs of known theorems.

The author lectured on quasi-conformal mappings at the University of Paris in the Fall of 1966, and wishes to acknowledge his debt to Professor C. Chevalley whose critical comments added to the clarity of this account.

## CONTENTS

## § 1. THE MOEBIUS GROUP

Definition. - Moebius $n$-space is the one-point compactification of euclidean $n$-space $\mathbf{R}^{n}$; it is denoted by $\mathbf{R}^{n} \cup \infty$.
$\mathrm{GM}(n)$, the Moebius group of Moebius $n$-space is the group of transformations generated by reflections in the spheres and $(n-1)$ planes of $\mathbf{R}^{n}$.
The composition of the reflection $\sigma_{\lambda}: p \rightarrow \frac{\lambda^{2} p}{|p|^{2}}$ of the sphere with center at the
origin and radius $\lambda$ and the reflection $\sigma_{1}$ is

$$
p \rightarrow \frac{\lambda^{2} p}{|p|^{2}} \rightarrow \frac{\lambda^{2} p}{|p|^{2}}\left|\frac{\lambda^{2} p}{|p|^{2}}\right|^{-2}=\lambda^{-2} p
$$

Similarly the composition of reflections in two parallel ( $n-1$ )-planes at a distance $d$ gives a translation by $2 d$ along the normal to the planes. Since the group generated by translations and stretchings permute the spheres of $\mathbf{R}^{n}$ transitively, and since $\sigma_{1}$ carries spheres through the origin into planes, we see easily that GM(n) is generated by $\sigma_{1}$ together with the subgroup of translations and stretchings.

Moebius space can be given a conformal structure invariant under GM(n) and compatible with the euclidean metric structure of $\mathbf{R}^{n}$.

Let $\mathrm{S}^{n}$ be a sphere in $\mathbf{R}^{n+1}$ i.e. $\left(\xi_{1}-a_{1}\right)^{2}+\ldots+\left(\xi_{n+1}-a_{n+1}\right)^{2}=r^{2}$, and let $p_{\infty}$ be a point of $\mathrm{S}^{n}$. Let F be an $n$-plane in $\mathbf{R}^{n+1}$ parallel to the tangent plane to $\mathrm{S}^{n}$ at $p_{\infty}$ and different from it. Stereographic projection from $p_{\infty}$ of $\mathrm{S}^{n}$ onto F is the mapping $\pi$ of $\mathrm{S}^{n}-p_{\infty}$ onto F given by projecting from $p_{\infty}$. The remarkable property of stereographic projection is that it preserves angles.

For let $L$ and $L^{\prime}$ be any two lines in $F$ through a point $x$. Let $N$ and $N^{\prime}$ denote the 2-planes determine by $p_{\infty} \cup \mathrm{L}$ and $p_{\infty} \cup L^{\prime}$ respectively. Then $N$ and $N^{\prime}$ cut $S^{n}$ in two circles C and $\mathrm{C}^{\prime}$ which pass through $p_{\infty}$ and $\pi^{-1}(x)$. The angle formed by C
and $\mathrm{C}^{\prime}$ at $\pi^{-1}(x)$ equals the angle formed by them at $p_{\infty}$. (This fact involves looking only at the three dimensional plane spanned by $p_{\infty} \cup L \cup L^{\prime}$ and its intersection with $\mathrm{S}^{n}$ which is an ordinary 2 -sphere.) On the other hand, the angle formed by C and $\mathrm{C}^{\prime}$ at $p_{\infty}$ equals the angle formed by $L, L^{\prime}$, since $F$ is parallel to the tangent plane to $S^{n}$ at $p_{\infty}$. Hence $\nless\left(\mathrm{L}, \mathrm{L}^{\prime}\right)$ equals the angle formed by C and $\mathrm{C}^{\prime}$ at $\pi^{-1}(x)$.

It is useful to have the analytic expression for stereographic projection. Let $\mathrm{S}^{n}$ be the $n$-sphere $\xi_{1}^{2}+\xi_{2}^{2}+\ldots+\xi_{n}^{2}+\left(\xi_{n+1}-\frac{1}{2}\right)^{2}=1 / 4$, let $\mathbf{R}^{n}$ be the plane $\xi_{n+1}=0$, let $\xi=\left(\xi_{1}, \ldots, \xi_{n+1}\right)$ be a point on $\mathrm{S}^{n}$, and let $x=\left(x_{1}, \ldots, x_{n}, o\right)=\pi(\xi), \pi$ denoting stereographic projection onto $\mathbf{R}^{n}$ from $p_{\infty}=(0, \ldots, \mathrm{I})$. Then

$$
\frac{\xi_{i}}{x_{i}}=\frac{\mathrm{I}-\xi_{n+1}}{\mathrm{I}} \quad \text { and } \quad \frac{\mathrm{I}-\xi_{n+1}}{\mathrm{I}}=\frac{\cos \theta}{\left(\mathrm{I}+|x|^{2}\right)^{1 / 2}}=\frac{\frac{\mathrm{I}}{\left(\mathrm{I}+|x|^{2}\right)^{1 / 2}}}{\left(\mathrm{I}+|x|^{2}\right)^{2 / 2}}=\frac{\mathrm{I}}{\mathrm{I}+|x|^{2}} .
$$

Thus
[ $\pi$ ]

$$
x_{i}=\frac{\xi_{i}}{\mathrm{I}-\xi_{n+1}} \quad(i=\mathrm{I}, \ldots, n)
$$

and
$\left[\pi^{-1}\right]$

$$
\xi_{i}=\frac{x_{i}}{\mathrm{I}+|x|^{2}}, \quad \xi_{n+1}=\frac{|x|^{2}}{\mathrm{I}+|x|^{2}} .
$$

Note. - Stereographic projection from ( $0,0, \ldots, 1$ ) of the unit sphere $\sum_{i=1}^{n+1} \xi_{i}^{2}=\mathrm{I}$ onto the plane $\xi_{n+1}=0$ is given by the same formula [ $\pi$ ]:


$$
\left\{\begin{aligned}
\xi_{i} & =\frac{\eta_{i}}{2} \quad(i=\mathrm{I}, \ldots, n) \\
\xi_{n+1} & =\frac{\eta_{n+1}+\mathrm{I}}{2}
\end{aligned}\right.
$$

[ $\pi]$

$$
x_{i}=\frac{\xi_{i}}{\mathrm{I}-\xi_{n+1}}=\frac{\frac{\eta_{i}}{2}}{\mathrm{I}-\frac{\eta_{n+1}+\mathrm{I}}{2}}=\frac{\eta_{i}}{\mathrm{I}-\eta_{n+1}} .
$$

But for $\pi^{-1}$ :


$$
\left\{\begin{array}{c}
\eta_{i}=2 \xi_{i}=\frac{2 x_{i}}{\mathrm{I}+|x|^{2}}, \quad i=\mathrm{I}, \ldots, n \\
\eta_{n+1}=2 \xi_{n+1}-\mathrm{I}=\frac{2|x|^{2}}{\mathrm{I}+|x|^{2}}-\mathrm{I}=\frac{|x|^{2}-\mathrm{I}}{|x|^{2}+\mathrm{I}}
\end{array}\right.
$$

If $\sigma_{1}$ is the reflection in the unit sphere of $\mathbf{R}^{n}$, then $\pi \sigma_{1} \pi^{-1}$ is the reflection of $S^{n}$ in its equatorial plane $\xi_{n+1}=\frac{1}{2}$ and is thus conformal, that is, it preserves angles. The stretching $h_{\lambda}: x \rightarrow \lambda x$ of $\mathbf{R}^{n}$ corresponds to the mapping of $\mathrm{S}^{n}$ given by $\quad \xi_{i} \rightarrow \frac{\lambda x_{i}}{1+\lambda^{2}|x|^{2}}$. The effect of this mapping for large $|x|$ is seen by expressing

$$
\frac{\lambda x_{i}}{\mathrm{I}+\lambda^{2}|x|^{2}}=\frac{\lambda x_{i}}{\mathrm{I}+|x|^{2}} \frac{\mathrm{I}+|x|^{2}}{\mathrm{I}+\lambda^{2}|x|^{2}}=\lambda \xi_{i} \cdot \frac{\mathrm{I}+|x|^{2}}{\mathrm{I}+\lambda^{2}|x|^{2}} .
$$

Thus at $p_{\infty}$, this mapping is approximately $\xi_{i} \rightarrow \lambda^{-1} \xi_{i}$. In particular it is analytic and conformal at $p_{\infty}$.

Regarding $\operatorname{GM}(n)$ as a transformation group on the $n$-sphere $\mathrm{S}^{n}$, one sees at once that it is the group generated by the group of rotations of $\mathrm{S}^{n}$, the reflection in some equatorial sphere, and the one parameter group corresponding to the stretchings $h_{\lambda}$.

Theorem ( $\mathbf{I} . \mathbf{1}$ ). - $\mathrm{GM}(n)$ is isomorphic as a transformation group to $\mathrm{O}(\mathrm{I}, n+\mathbf{1}) /( \pm \mathbf{1})$, the orthogonal group of the quadratic form $y_{0}^{2}-y_{1}^{2} \ldots-y_{n+1}^{2}$, operating on the projective variety $y_{0}^{2}-y_{1}^{2}-\ldots-y_{n+1}^{2}=\mathrm{o} . \quad$ Moreover, if we set $\eta_{i}=\frac{y_{i}}{y_{0}}(i=\mathrm{I}, \ldots, n+\mathrm{I})$, then $\mathrm{O}(\mathrm{I}, n+\mathrm{I}) /( \pm \mathrm{I})$
is the Moebius group of $n_{1}^{2}+\ldots+\eta^{2}=1$. is the Moebius group of $\eta_{1}^{2}+\ldots+\eta_{n+1}^{2}=\mathrm{I}$.

Proof. - Introducing non-homogeneous coordinates, we can identify $\mathrm{O}(\mathrm{I}, n+\mathrm{I})^{\prime}( \pm \mathrm{r})$ with a group of transformations of the $n$-sphere $\mathrm{S}^{n}: \eta_{1}^{2}+\ldots \eta_{n+1}^{2}=1 . \quad \mathrm{O}(\mathrm{I}, n+1)$ clearly contains the group $\mathrm{O}(n+\mathrm{I})$, the orthogonal group of an $\mathrm{S}^{n}$; also $\mathrm{O}(\mathrm{I}, n+\mathrm{I})$ contains the reflection $\eta_{1} \rightarrow-\eta_{1}$. To prove that $\mathrm{O}(\mathrm{I}, n+\mathrm{I}) /( \pm \mathrm{I})$ contains $\mathrm{GM}(n)$, it suffices only to verify that it contains the one parameter group corresponding to stretchings of $\mathbf{R}^{n}$.

Let $h_{\lambda}^{\prime}$ denote the element of $\mathrm{GL}(n+2, \mathbf{R})$ given by

$$
\begin{aligned}
y_{0}+y_{n+1} & \rightarrow \lambda\left(y_{0}+y_{n+1}\right) \\
y_{0}-y_{n+1} & \rightarrow \lambda^{-1}\left(y_{0}-y_{n+1}\right) \\
y_{i} & \rightarrow y_{i} \quad(i=\mathrm{I}, \ldots, n) .
\end{aligned}
$$

Taking stereographic projection from ( $0,0, \ldots, \mathrm{I}$ ) onto the plane $\eta_{n+1}=0$, we get

$$
x_{i}^{\prime}=\frac{\eta_{i}^{\prime}}{\mathrm{I}-\eta_{n+1}^{\prime}}=\frac{y_{i}^{\prime}}{y_{0}^{\prime}-y_{n+1}^{\prime}}=\lambda \cdot \frac{y_{i}}{y_{0}-y_{n+1}}=\lambda x_{i} \quad(i=\mathrm{I}, \ldots, n)
$$

which is the stretching $h_{\lambda}$ on $\mathbf{R}^{n}$.
Theorem (1.2). - The subgroup $\mathrm{G}^{\prime}$ of $\mathrm{GM}(n)$ which stabilizes the hemisphere $\mathrm{S}_{-}: \eta_{n+1}<0$ is isomorphic to $\mathrm{GM}(n-1)$ under the restriction homomorphism into its action on the equatorial $(n-1)$-sphere: $\quad \eta_{n+1}=0$. Moreover $\mathrm{G}^{\prime}$ operates transitively on $\mathrm{S}_{-}$and keeps invariant a positive definite quadratic differential form ds ${ }^{2}$. Under stereographic projection from ( $0,0, \ldots, \mathrm{I}$ ), $\mathrm{S}_{-}$maps onto the unit ball $x_{1}^{2}+\ldots+x_{n}^{2}<1$ and its invariant metric $d s^{2}$ becomes $\frac{d x_{1}^{2}+\ldots+d x_{n}^{2}}{\left(1-|x|^{2}\right)^{2}}$
(up to a constant factor). (up to a constant factor).

Proof. - We continue the notations employed above. Then $\mathrm{G}^{\prime}$ can be identified modulo a factor $\pm \mathrm{I}$, with the subgroup of $\mathrm{O}(\mathrm{I}, n+\mathrm{I})$ which keeps invariant the halfspace $y_{n+1}<0$ in the $n+2$ dimensional cartesian space $\left(y_{0}, \ldots, y_{n+1}\right)$. It follows that any element $g$ in this subgroup $\mathrm{G}^{\prime}$ sends $y_{n+1}$ into $c y_{n+1}, c>0$. Since $g$ keeps invariant the plane $y_{n+1}=0$ as well as its orthogonal complement, it follows from the invariance of $y_{0}^{2}-\ldots-y_{n+1}^{2}$ that $c=\mathrm{I}$ and $\mathrm{G}^{\prime}$ preserves $y_{0}^{2}-y_{1}^{2}-\ldots-y_{n}^{2}$; that is, $\mathrm{G}^{\prime}$ can be identified with $\mathrm{O}(\mathrm{I}, n) /( \pm \mathrm{I})$; by the preceding theorem, $\mathrm{G}^{\prime} \approx \mathrm{GM}(n-\mathrm{I})$.

We next define the $\mathrm{G}^{\prime}$-invariant metric on $\mathrm{S}_{-}$. In $\mathbf{R}^{n+2}$, let A denote one of the two connected components of the $n$-dimensional surface

$$
\left\{\begin{array}{l}
y_{0}^{2}-y_{1}^{2}-\ldots-y_{n+1}^{2}=\mathbf{o} \\
y_{n+1}=-1
\end{array}\right.
$$

say the half defined by $y_{0}>0$. Then A is a cross-section for the fiber map of the half-cone

$$
\left\{\begin{array}{l}
y_{0}^{2}-y_{1}^{2}-\ldots-y_{n+1}^{2}=0 \\
y_{n+1}<0
\end{array}\right.
$$

into the corresponding projective variety, and indeed A maps onto the hemisphere $\mathrm{S}_{-}$:

$$
\left\{\begin{array}{l}
\eta_{1}^{2}+\ldots+\eta_{n+1}^{2}=\mathbf{x} \\
\eta_{n+1}<0
\end{array} .\right.
$$

The differential form $d s_{\mathrm{L}}^{2}=d y_{0}^{2}-d y_{1}^{2}-\ldots-d y_{n}^{2}$ is clearly invariant under $\mathrm{G}^{\prime}=\mathrm{O}(\mathrm{I}, n) \bmod ( \pm \mathrm{I})$. We compute the form induced by $d s_{\mathrm{L}}^{2}$ on $\mathrm{S}_{-}$:

From $y_{i}=n_{i} y_{0}(i=1, \ldots, n)$, and $y_{0}^{2}-\sum_{i=1}^{n} y_{i}^{2}=\mathrm{r}$, we get

$$
\begin{aligned}
& y_{0}^{2}\left(\mathrm{I}-\sum_{i=1}^{n} \eta_{i}^{2}\right)=\mathrm{I} \\
& y_{0}=\left(\mathrm{I}-\sum_{i=1}^{n} \eta_{i}^{2}\right)^{-1 / 2}
\end{aligned}
$$

$$
\begin{gathered}
d y_{0}=\frac{\sum_{i} \eta_{i} d \eta_{i}}{\left(\mathrm{I}-\sum_{i} \eta_{i}^{2}\right)^{3 / 2}} \\
d y_{i}=y_{0} d \eta_{i}+\eta_{i} d y_{0}=\frac{d \eta_{i}}{\left(\mathrm{I}-\sum_{j} \eta_{j}^{2}\right)^{1 / 2}}+\frac{\eta_{i} \sum_{j} \eta_{j} d \eta_{j}}{\left(\mathrm{I}-\sum_{j} \eta_{j}^{2}\right)^{3 / 2}} \\
=\left(\mathrm{I}-\sum_{j} \eta_{j}^{2}\right)^{-3 / 2}\left(\eta_{i} \sum_{j} \eta_{j} d \eta_{j}+\left(\mathrm{I}-\sum_{j} \eta_{j}^{2}\right) d \eta_{i}\right) \\
d y_{0}^{2}-\sum_{i=1}^{n} d y_{i}^{2}=\left(\mathrm{I}-\sum_{j} \eta_{j}^{2}\right)^{-3}\left[\left(\sum_{j} \eta_{j} d \eta_{j}\right)^{2}-\sum_{i}\left(\eta_{i} \sum_{j} \eta_{j} d \eta_{j}+\left(\mathrm{I}-\sum_{\eta_{j}^{2}}^{2}\right) d \eta_{i}\right)^{2}\right] \\
=\left(\mathrm{I}-\sum_{j} \eta_{j}^{2}\right)^{-3}\left[\left(\sum_{\eta_{j}} d \eta_{j}\right)^{2}-\left(\sum_{i} \eta_{i}^{2}\right)\left(\sum_{j} \eta_{j} d \eta_{j}\right)^{2}-\left(\mathrm{I}-\sum_{j} \eta_{j}^{2}\right)^{2}\left(\sum_{i} d \eta_{i}^{2}\right)\right. \\
\left.=-2\left(\sum_{i} \eta_{i} d \eta_{i}\right)\left(\sum_{j} \eta_{j} d \eta_{j}\right)+2\left(\sum_{i} \eta_{i}^{2}\right)\left(\sum_{i} \eta_{i} d \eta_{i}\right)^{2}\right] \\
=-\left(\mathrm{I}-\sum_{j} \eta_{j}^{2}\right)^{-3}\left[( \sum _ { i } \eta _ { i } d \eta _ { i } ) ^ { 2 } \left(-\mathrm{I}+\sum_{j}^{2} \eta^{-1}\left[\left(\mathrm{I}-\sum_{j=1}^{2}\right)-\sum_{j}^{n} d \eta_{j}^{2}\left(\mathrm{I}-\sum_{j}^{2} \eta_{j}^{2}\right)^{2}\right]\right.\right. \\
-1 \\
i=1 \\
n \\
\left.\left.\eta_{i} d \eta_{i}\right)^{2}+\sum_{i=1}^{n} d \eta_{i}^{2}\right]
\end{gathered}
$$

The term in the [ ] is exactly $\sum_{i=1}^{n+1} d \eta_{i}^{2}$ on the sphere $\eta_{1}^{2}+\ldots+\eta_{n+1}^{2}=1$ since $\eta_{n+1}=\left(\mathrm{I}-\sum_{i=1}^{n} \eta_{i}^{2}\right)^{1 / 2} . \quad$ Hence $d s_{\mathrm{L}}^{2}$ becomes $-\eta_{n+1}^{-2}\left(d \eta_{1}^{2}+\ldots+d \eta_{n+1}^{2}\right)$.

Now

$$
\begin{aligned}
\eta_{i} & =\frac{2 x_{i}}{\mathrm{I}+|x|^{2}}, \quad i=\mathrm{I}, \ldots, n \\
\eta_{n+1} & =\frac{|x|^{2}-\mathrm{I}}{|x|^{2}+\mathrm{I}}
\end{aligned}
$$

Hence

$$
\begin{aligned}
& \begin{aligned}
\sum_{i=1}^{n+1} d \eta_{i}^{2} & =4 \sum_{i=1}^{n}\left(\frac{\left(\mathrm{I}+|x|^{2}\right) d x_{i}-2 x_{i} \sum_{j} x_{j} d x_{j}}{\left(\mathrm{I}+|x|^{2}\right)^{2}}\right)^{2}+\left(\frac{\left(\mathrm{I}+|x|^{2}\right) 2 \sum_{j} x_{j} d x_{j}-\left(|x|^{2}-\mathrm{I}\right) 2 \sum_{j} x_{j} d x_{j}}{\left(\mathrm{I}+|x|^{2}\right)^{2}}\right)^{2} \\
& =4\left(\mathrm{I}+|x|^{2}\right)^{-4}\left(\sum_{i=1}^{n}\left[\left(\mathrm{I}+|x|^{2}\right)^{2} d x_{i}^{2}-4\left(\mathrm{I}+|x|^{2}\right)\left(\sum_{x_{j}} d x_{j}\right) x_{i} d x_{i}+4\left(\sum_{j} x_{j} d x_{j}\right)^{2} x_{i}^{2}\right)\right] \\
& +\mathrm{I} 6\left(\mathrm{I}+|x|^{2}\right)^{-4}\left(\sum_{j},\right. \\
& =4\left(\mathrm{I}+|x|^{2}\right)^{-2} \sum_{i=1}^{n} d x_{i}^{2}
\end{aligned} \\
& \text { and } \eta_{n+2}^{-2}\left(\sum_{i=1}^{n+1} d \eta_{i}^{2}\right)=\frac{4}{\left(\mathrm{I}-|x|^{2}\right)^{2}} \sum_{i=1}^{n} d x_{i}^{2} .
\end{aligned}
$$

Remark 1. - We will prove later that $\mathrm{GM}(n)$ is the group of all conformal mappings of $S^{n}$ to $S^{n}$ for $n>1$.

Remark 2. - The hemisphere $\eta_{n+1}<0$ with metric $\eta_{n+1}^{-2} \sum_{i=1}^{n+1} d \eta_{i}^{2}$ or equivalently the unit $n$-ball $|x|<\mathrm{I}$ with metric $\left(\mathrm{I}-\left|x^{2}\right|\right)^{-2} \sum_{i=1}^{n} d x_{i}^{2}$ is a Riemannian space in which the isotropy subgroup at a point is $\mathrm{O}(n)$. Hence these spaces have constant curvature.

Definition. - n-dimensional hyperbolic space is the Riemannian space described in Remark 2.

## § 2. CONFORMAL CAPACITY OF A SHELL

Definition. - A shell in Moebius n-space is a connected open set D whose complement consists of two connected components $\mathrm{C}_{0}$ and $\mathrm{C}_{1}$. The subsets $\Delta_{0}=\mathrm{C}_{0} \cap \overline{\mathrm{D}}$ and $\Delta_{1}=\mathrm{C}_{1} \cap \overline{\mathrm{D}}$ are called the boundaries of D .

A shell not containing the point $\infty$ is called a shell in $\mathbf{R}^{n}$, the component $\mathrm{G}_{i}$ of its complement which contains $\infty$ is the unbounded component and $\Delta_{i}$ is called the boundary of the unbounded component of the complement.

Let $D$ be a shell in the Moebius space $S^{n}$. Injection of the pair ( $\overline{\mathrm{D}}, \overline{\mathrm{D}}-\mathrm{D}$ ) into ( $\mathrm{S}^{n}, \mathrm{~S}^{n}-\mathrm{D}$ ) gives rise to the following commutative diagram of exact integral cohomology sequences based on compact supports:


It follows that $H_{c}^{1}(\mathbf{D})=\mathbf{Z}$, the ring of integers, that

$$
\mathrm{H}^{0}\left(\Delta_{0} \cup \Delta_{1}\right)=\mathrm{H}^{0}(\overline{\mathrm{D}}-\mathrm{D}) \approx \mathrm{H}^{0}\left(\mathrm{~S}^{n}-\mathrm{D}\right) \approx \mathbf{Z}+\mathbf{Z}
$$

and thus each $\Delta_{i}$ is connected.
Note. - From the fact that $\mathrm{H}_{c}^{1}(\mathbf{D})=\mathbf{Z}$ and the invariance of domain theorem, it follows that any homeomorphism of a shell D into $\mathrm{S}^{n}$ has a shell for its image. The homeomorphism need not, of course, extend to a homeomorphism on the boundaries.

We adopt the following conventions for notation. If $f$ is a $\mathrm{C}^{1}$ mapping of manifolds, $\dot{f_{p}}$ denotes the differential of $f$ at the point $p$, or merely $\dot{f}$ when the intended point $p$ is clear from the context. In the special case of $\mathrm{C}^{1}$ maps of $\mathbf{R}^{n}$ to $\mathbf{R}^{m}$, we identify $\dot{f_{p}}$ with a map of $\mathbf{R}^{n}$ to $\mathbf{R}^{m}$ in the standard way; that is, by identifying the tangent space at any point of Cartesian space $\mathbf{R}^{n}$ with $\mathbf{R}^{n}$ itself. We denote by $\langle\mathrm{X}, \mathrm{Y}\rangle$ the standard inner product on either $\mathbf{R}^{n}$ or its dual space. We write $|\mathrm{X}|=\langle\mathrm{X}, \mathrm{X}\rangle^{1 / 2}$.

Definition. - Let D be a shell in Moebius n-space and let $\Delta_{0}, \Delta_{1}$ denote its boundaries. The conformal capacity $\mathscr{C}(\mathrm{D})$ is defined as

$$
\operatorname{Inf}_{u} \int_{\mathrm{D}}|\nabla u|^{n} d \mathrm{D}
$$

where $u$ is a continuous real valued function on $\overline{\mathrm{D}}$ which is of class $\mathrm{C}^{1}$ on D and takes the constant value o on $\Delta_{0}$ and the constant value I on $\Delta_{1}$; $\nabla u$ denotes the gradient of $u$ with respect to the standard metric of $\mathbf{R}^{n}$, and $d \mathrm{D}$ denotes the standard measure on $\mathbf{R}^{n}$.

Note. - The definition of $\mathscr{C}(\mathrm{D})$ does not depend on designation of the boundaries since $|\nabla u|=|\nabla(\mathrm{r}-u)|$.

Given a diffeomorphism of shells

$$
f: \mathbf{D} \rightarrow \mathbf{D}^{\prime}
$$

and given $u$, a $\mathrm{C}^{1}$ function on D taking boundary values o and I on $\Delta_{0}$ and $\Delta_{1}$, we define the function $u^{\prime}$ on $\mathrm{D}^{\prime}$ so that

$$
u=u^{\prime} \circ f
$$

It is easily seen that $u^{\prime}$ is a $\mathbf{C}^{1}$ function on the shell $\mathrm{D}^{\prime}$ taking the value o on one boundary and the value I on the other. Writing $y=f(x)$, we have

$$
\nabla u=\left(\frac{\partial u}{\partial x_{i}}\right), \quad \frac{\partial u}{\partial x_{i}}=\sum_{j} \frac{\partial u^{\prime}}{\partial y_{j}} \frac{\partial y_{j}}{\partial x_{i}},
$$

so that $\nabla u={ }^{t} \dot{f}\left(\nabla u^{\prime}\right)$, where ${ }^{t} \dot{f}$, the transpose of $\dot{f}$, sends covectors at $f(p)$ to covectors at $p$. Hence

$$
\begin{aligned}
\int_{\mathrm{D}^{\prime}}\left|\nabla u^{\prime}\right|^{n} d \mathrm{D}^{\prime} & =\int_{\mathrm{D}}\left(\left|\nabla u^{\prime}\right|^{n} \circ f\right)|\operatorname{det} \dot{f}| d \mathrm{D} \\
& =\int_{\mathrm{D}}\left(\left|\dot{f}^{-1}(\nabla u)\right|^{n} \circ f\right)|\operatorname{det} \dot{f}| d \mathrm{D} .
\end{aligned}
$$

Suppose now that $f$ is conformal, that is, for all points $p \in \mathrm{D}$ and for all non-zero tangent vectors X and Y at $p$,

$$
\frac{\left\langle\dot{f_{p}} \mathrm{X}, \dot{f_{p}} \mathrm{Y}\right\rangle}{\left|\dot{f_{p}} \mathrm{X}\right| \cdot\left|\dot{f_{p}} \mathrm{Y}\right|}=\frac{\langle\mathrm{X}, \mathrm{Y}\rangle}{|\mathrm{X}| \cdot|\mathrm{Y}|}
$$

For any two orthogonal unit vectors $\mathrm{X}_{i}, \mathrm{X}_{j}$ we deduce from $\left\langle\mathrm{X}_{i}+\mathrm{X}_{j}, \mathrm{X}_{i}-\mathrm{X}_{j}\right\rangle=0$ that

$$
o=\left\langle\dot{f_{p}}\left(\mathrm{X}_{i}+\mathrm{X}_{j}\right), \dot{f_{p}}\left(\mathrm{X}_{i}-\mathrm{X}_{\mathrm{j}}\right)\right\rangle=\left|\dot{f_{p}}\left(\mathrm{X}_{i}\right)\right|^{2}-\left|\dot{f_{p}}\left(\mathrm{X}_{j}\right)\right|^{2}
$$

It follows that $\dot{f_{p}}$ maps the unit ball in the tangent space at $p$ to a ball of radius $\lambda(p)=\left|\dot{f_{p}}\left(\mathrm{X}_{i}\right)\right|$; equivalently, ${ }^{t} \dot{f}_{p} \dot{f}_{p}=\lambda^{2}(p)$.identity.

Hence if $f: \mathrm{D} \rightarrow \mathrm{D}^{\prime}$ is a conformal mapping, we have

$$
\left|\operatorname{det} \dot{f}_{p}\right|^{2}=\operatorname{det}\left(\dot{f}_{p} \cdot \dot{f_{p}}\right)=\lambda^{2 n}(p)
$$

and

$$
|\nabla u|^{2}=\left\langle^{\imath} \dot{f}\left(\nabla u^{\prime}\right),{ }^{\imath} \dot{f}\left(\nabla u^{\prime}\right)\right\rangle=\left\langle\dot{f}^{.} \cdot \dot{f}\left(\nabla u^{\prime}\right), \nabla u^{\prime}\right\rangle .
$$

But $\dot{f}^{t} \cdot \dot{f}={ }^{t} \dot{f}^{-1}\left({ }^{t} \dot{f} \cdot \dot{f}\right)^{t} \dot{f}={ }^{t} \dot{f}^{-1} \cdot \lambda^{2} \cdot{ }^{t} \dot{f}=\lambda^{2}$. Hence $|\nabla u|^{2}(p)=\lambda^{2}(p)\left|\nabla u^{\prime}\right|^{2}(f(p))$, that is

$$
|\nabla u|^{n}=|\operatorname{det} \dot{f}|\left(\left|\nabla u^{\prime}\right|^{n} \circ f\right)
$$

Hence

$$
\int_{\mathrm{D}^{\prime}}\left|\nabla u^{\prime}\right|^{n} d \mathrm{D}^{\prime}=\int_{\mathrm{D}}|\nabla u|^{n} d \mathrm{D}
$$

if $f$ is conformal.
Note. - Conformal capacity in dimension 2 is a classical notion identical with Newtonian capacity. For $n>2$, it was introduced by C. Loewner (Journal of Mathematics and Mechanics, vol. 8, 1959) who proved that $\mathscr{C}(\mathrm{D})>0$ if and only if neither boundary of D degenerates to a point. We shall prove this later.

We shall first present some technical lemmas about the functions $u$ that define $\mathscr{C}(\mathrm{D})$; these are needed to establish the important fact that $\mathscr{C}(\mathrm{D})$ varies continuously with D .

## § 3. ADMISSIBLE FUNCTIONS

A continuous function $f$ on the interval $a \leqslant x \leqslant b$ is called absolutely continuous if its derivative $\frac{d f}{d x}$ exists almost everywhere and is integrable and $\int_{x_{0}}^{x_{1}} \frac{d f}{d x} d x=f\left(x_{1}\right)-f\left(x_{0}\right)$ for all $a \leqslant x_{0}, x_{1} \leqslant b$. It should be pointed out that if $\frac{d f}{d x}$ exists almost everywhere, then it is measurable (Saks, p. II2). This condition is equivalent to the assertion: Given an $\varepsilon>0$, there is a $\delta>0$ such that for any union of intervals $A$ of measure less than $\delta$, we have $\mathrm{F}(\mathrm{A})<\varepsilon$, where $\mathbf{F}$ is the set function formed from $f$ by the rule:

$$
\mathrm{F}\left[x_{0}, x_{1}\right]=f\left(x_{1}\right)-f\left(x_{0}\right)
$$

A continuous function $f$ defined on a domain $\mathrm{D} \subset \mathbf{R}^{n}$ is called ACL in D if in any closed ball lying in D it is absolutely continuous on almost all lines parallel to the coordinate axes. (If in addition we assume that $\int_{\mathrm{D}}|\nabla f| d \mathrm{D}<\infty$, the condition is invariant under diffeomorphisms, but we shall not need this fact.)

The partial derivates of an ACL function exist almost everywhere, as is seen by applying Fubini's theorem to the set on which the derivatives do not exist.

With ACL functions having integrable derivatives, one can differentiate under an integral sign.

Lemma (3.1). - Let D be a domain in $\mathbf{R}^{n}$. Let $f$ be continuous on $[a, b] \times \mathrm{D}$, and absolutely continuous in $x$ for almost all $y$ in D . Assume $\frac{\partial f}{\partial x}(x, y)$ is integrable on $[a, b] \times \mathrm{D}$. Then $\frac{d}{d x} \int_{\mathrm{D}} f(x, y) d y=\int_{\mathrm{D}} \frac{\partial f}{\partial x}(x, y) d y$ for almost all $x$ in $[a, b]$.

Proof. - By hypothesis ACL, $\left(\frac{\partial f}{\partial x}(x, y)\right.$ is integrable in $x$ for almost all $y \in \mathrm{D}$, and $)$ $\int_{x_{0}}^{x} \frac{\partial f}{\partial x}(x, y) d x=f(x, y)-f\left(x_{0}, y\right)$ for almost all $y$. Since by hypothesis $\frac{\partial f}{\partial x}(x, y)$ is integrable
on $[a, b] \times \mathrm{D}$, we find by Fubini's theorem (Saks, p. 8I) that $g(x)=\int_{\mathrm{D}} \frac{\partial f}{\partial x}(x, y) d y$ exists for almost all $x$, is integrable on $[a, b]$ and that

$$
\int_{x_{0}}^{x} g(x) d x=\int_{\mathrm{D}}\left(\int_{x_{0}}^{x} \frac{\partial f}{\partial x}(x, y) d x\right) d y=\int_{\mathrm{D}}\left(f(x, y)-f\left(x_{0}, y\right) d y\right.
$$

Since (Saks, p. II7), $g(x)=\frac{d}{d x} \int_{x_{0}}^{x} g(x) d x$ for almost all $x$, we get

$$
\frac{d}{d x} \int_{\mathrm{D}} f(x, y) d y=\int_{\mathrm{D}} \frac{\partial f}{\partial x}(x, y) d y
$$

for almost all $x$ in $[a, b]$.
The next lemma collects some elementary observations that will be used repeatedly.
Lemma (3.2). - Let v be a continuous, ACL function in an open set $\mathrm{R} \subset \mathbf{R}^{n}$. Let $\mathbf{R}^{\prime}$ be an open subset with compact closure and $\overline{\mathbf{R}}^{\prime} \subset \mathbf{R}$. Let U be the ball $|y|<\varepsilon$ in $\mathbf{R}^{n}$, and suppose that $2 \varepsilon$ does not exceed the distance from $\mathbf{R}^{\prime}$ to the complement of $\mathbf{R}$. Assume that $|\nabla v|$ is integrable on R . Set $w(x)=\frac{\mathrm{I}}{m(\mathrm{U})} \int_{\mathrm{U}} v(x+y) d y$ for $x+\mathrm{U} \subset \mathrm{R}$, where $m(\mathrm{U})$ denotes the $n$-dimensional measure of U .

Then

1) $w$ is of class $\mathrm{C}^{1}$ on $\mathrm{R}^{\prime}$.
2) $\lim _{\varepsilon \rightarrow 0} w=v$ uniformly on compact subsets of $\mathbf{R}$.
3) $\nabla w(x)=\frac{\mathrm{I}}{m(\mathrm{U})} \int_{\mathrm{U}} \nabla v(x+y) d y$ for almost all $x$.
4) $\int_{\mathbf{R}^{\prime}}|\nabla w(x)|^{p} d x \leqslant \int_{\mathbf{R}}|\nabla v(x)|^{p} d x \quad p \geqslant 1$.

Proof. - Assertions 1) and 2) are well-known facts about continuous functions. Assertion 3) follows directly from Lemma (3.1).

We have $|\nabla w(x)| \leqslant \frac{\mathrm{I}}{m(\mathrm{U})} \int_{U}|\nabla v(x+y)| d y$. Hence

$$
\begin{aligned}
\left(\int_{\mathrm{R}^{\prime}}|\nabla w(x)|^{p} d x\right)^{\frac{1}{p}} & \leqslant\left(\int_{\mathrm{R}^{\prime}}\left(\frac{\mathrm{I}}{m(\mathrm{U})} \int_{\mathrm{U}}|\nabla v(x+y)| d y\right)^{p} d x\right)^{\frac{1}{p}} \\
& \leqslant \frac{\mathrm{I}}{m(\mathrm{U})} \int_{\mathrm{U}}\left(\int_{\mathrm{R}^{\prime}}|\nabla v(x+y)|^{p} d x\right)^{\frac{1}{p}} d y
\end{aligned}
$$

by the integral form of Minkowski's inequality $\left\|\sum_{i} f_{i}\right\|_{p} \leqslant \sum_{i}\left\|f_{i}\right\|_{p}$.
Now $\int_{\mathbf{R}^{\prime}}|\nabla v(x+y)|^{p} d x \leqslant \int_{\mathrm{R}}|\nabla v(x)|^{p} d x$; hence we get

$$
\left(\int_{\mathrm{R}^{\prime}}|\nabla w(x)|^{p} d x\right)^{\frac{1}{p}} \leqslant\left(\int_{\mathrm{R}}|\nabla v(x)|^{p} d x\right)^{\frac{1}{p}} \frac{\mathrm{I}}{m(\mathrm{U})} \int_{\mathrm{U}} d y
$$

from which assertion 4) follows.
Let D be a shell in $\mathbf{R}^{n}$, let $\mathrm{C}_{1}$ be the unbounded component of its complement in Moebius space, and let $\mathrm{C}_{0}$ be the bounded component.

A real valued function $u$ defined on $\overline{\mathrm{D}}$ is called admissible if it is continuous, ACL in D , takes the value o on $\mathrm{C}_{0} \cap \overline{\mathrm{D}}$ and the value I on $\mathrm{C}_{1} \cap \overline{\mathrm{D}}$ and if moreover $|\nabla u|^{n}$ is integrable on D . The function $u$ is called smoothly admissible if in addition it is of class $\mathrm{C}^{1}$ on D and its gradient $\nabla u$ has compact support in D .

The shell $D$ has an admissible function. To see this, cover $\mathrm{C}_{1}$ by a finite set of closed balls in Moebius space whose union $\mathrm{C}_{1}^{\prime}$ does not meet $\mathrm{C}_{0}$ (a closed ball with center at $\infty$ being $\{y ;|y| \geqslant r\}$, and then cover $\mathrm{C}_{0}$ with a finite number of closed balls whose union $\mathrm{C}_{0}^{\prime}$ does not meet $\mathrm{C}_{1}^{\prime}$. The functions $d\left(x, \mathrm{C}_{0}^{\prime}\right)$ and $d\left(x, \mathrm{C}_{1}^{\prime}\right)$ are then easily seen to be ACL. In fact, for any set C , the function $f(x)=d(x, \mathrm{C})$ satisfies $|f(x+\Delta x)-f(x)| \leqslant|\Delta x|$ by the triangle inequality, and from this it is seen that $f(x)$ is absolutely continuous on all lines. Hence $u(x)=\frac{d\left(x, \mathrm{C}_{0}^{\prime}\right)}{d\left(x, \mathrm{C}_{0}^{\prime}\right)+d\left(x, \mathrm{C}_{1}^{\prime}\right)}$ is ACL and its gradient is a.e. bounded and of compact support. Thus $\int_{\mathrm{D}}|\nabla u|^{n} d \mathrm{D}<\infty$ and $u$ is admissible.

Lemma (3.3). - The conformal capacity of the shell D may be defined by

$$
\operatorname{Inf}_{u} \int_{\mathrm{D}}|\nabla u|^{n} d \mathrm{D}
$$

where $u$ ranges over the class of admissible functions or over the class of smoothly admissible functions.
Proof. - Let $u$ be an admissible function.
Fix $a$ : $0<a<\frac{1}{2}$. Set

$$
v(x)= \begin{cases}0 & \text { if } u(x) \leqslant a \\ \frac{u(x)-a}{\mathrm{I}-2 a} & \text { if } a<u(x)<\mathrm{I}-a \\ \mathrm{I} & \text { if } \mathrm{I}-a \leqslant u(x)\end{cases}
$$

and extend $v$ to be o on $\mathrm{C}_{0}$ and I on $\mathrm{C}_{1}$. Then $\nabla v$ has compact support K , which lies in D at a positive distance, say $2 \varepsilon$, from the complement of D . Let $\mathrm{U}=\{y ;|y|<\varepsilon\}$. Set

$$
w(x)=\frac{\mathrm{I}}{m(\mathrm{U})} \int_{\mathrm{U}} v(x+y) d y
$$

Then $w=0$ on $\mathrm{C}_{0}$ and $w=1$ on $\mathrm{C}_{1}$. Since $\nabla v$ has compact support K , we get by Hölder's inequality

$$
\begin{aligned}
\int_{\mathrm{D}}|\nabla v| d \mathrm{D} & =\int_{\mathrm{K}}|\nabla v| \cdot \mathrm{I} d \mathrm{D} \leqslant\left(\int_{\mathrm{K}}|\nabla v|^{n} d \mathrm{D}\right)^{\frac{1}{n}}\left(\int_{\mathrm{K}} \mathrm{I} \cdot d \mathrm{D}\right)^{\frac{n-1}{n}} \\
& \leqslant(\mathrm{I}-2 a)^{-1}\left(\int_{\mathrm{D}}|\nabla u|^{n} d \mathrm{D}\right)^{\frac{1}{n}} m(\mathrm{~K})<\infty
\end{aligned}
$$

Thus $|\nabla v|$ is integrable on D. Hence for almost all $x$,

$$
\nabla w(x)=\frac{\mathrm{I}}{m(\mathrm{U})} \int_{\mathrm{U}} \nabla v(x+y) d y
$$

by Lemma 3.2, and moreover $w$ is of class $\mathrm{C}^{1}$ on $\mathrm{K}+\overline{\mathrm{U}}$, which contains the support of $\nabla w$. Thus $w$ is a smoothly admissible function on D. Hence

$$
\mathscr{C}(\mathrm{D}) \leqslant \int_{\mathrm{D}}|\nabla w|^{n} d x \leqslant \int_{\mathrm{D}}|\nabla v|^{n} d x
$$

by assertion 4) of Lemma (3.2), and on the other hand

$$
\int_{D}|\nabla v|^{n} d x \leqslant \int_{D}|\nabla u|^{n}(\mathrm{I}-2 a)^{-n} d x
$$

Hence

$$
\mathscr{C}(\mathrm{D}) \leqslant \int_{\mathrm{D}}|\nabla w|^{n} d x \leqslant(\mathrm{I}-2 a)^{-n} \int_{\mathrm{D}}|\nabla u|^{n} d x
$$

Letting $a \rightarrow 0$, we find

$$
\mathscr{C}(\mathbf{D}) \leqslant \operatorname{Inf}_{u} \int_{D}|\nabla u|^{n} d x
$$

as $n$ varies over either the class of admissible functions or the class of smoothly admissible function.

Note. - We can now assert that $\mathscr{C}(\mathrm{D})<\infty$ for any shell D in Moebius space. For, performing a Moebius transformation, we can assume $\mathrm{D} \subset \mathbf{R}^{n}$. Then we can find an admissible function for $D$ as remarked above. Hence $\mathscr{C}(D)<\infty$ by Lemma (3.3).

Lemma (3.4). - Let D be a shell in Moebius space, let $\mathrm{C}_{\mathbf{0}}$ and $\mathrm{C}_{\mathbf{1}}$ denote the components of its complement. Let $u$ be an admissible function on D . Extend $u$ to a function on Moebius space by assigning it the constant values o and I on $\mathrm{C}_{0}$ and $\mathrm{C}_{1}$. Then $u$ is continuous and ACL everywhere on $\mathbf{R}^{n}$, and $|\nabla u|=0$ on $\mathrm{C}_{0} \cup \mathrm{C}_{1}$ a.e.

Proof. - Clearly $u$ is continuous.
Let X denote a variable line parallel to the $x_{1}$-axis. Then

$$
\int_{\text {Proj } \mathrm{D}}\left(\int_{\mathrm{X} \cap \mathrm{D}}|\nabla u|^{n} d x_{1}\right) d x_{2} \ldots d x_{n}=\int_{\mathrm{D}}|\nabla u|^{n} d x<\infty
$$

Hence

1) $\int_{\mathrm{X} \cap \mathrm{D}}|\nabla u|^{n} d x_{1}<\infty$ for almost all X . Morever, since $u$ is ACL in D ,
2) $u$ is absolutely continuous on any compact interval in $\mathrm{X} \cap \mathrm{D}$ for almost all X . Choose X satisfying 1) and 2).

Then given any interval $\mathrm{I}=[p, q] \subset \mathrm{X}$, we show

$$
|u(p)-u(q)| \leqslant \int_{\mathrm{I} \cap \mathrm{D}}|\nabla u| d x_{1} .
$$

Clearly no generality is lost in assuming $p \in \overline{\mathrm{D}}$ and $q \in \overline{\mathrm{D}}$. Let $p^{\prime}$ be the nearest point of $\left(\mathrm{C}_{0} \cup \mathrm{C}_{1}\right) \cap \mathrm{I}$ to $p$ and let $q^{\prime}$ be the nearest point of $\left(\mathrm{C}_{0} \cup \mathrm{C}_{1}\right) \cap \mathrm{I}$ to $q$. Then

$$
|u(p)-u(q)| \leqslant\left|u(p)-u\left(p^{\prime}\right)\right|+\left|u\left(p^{\prime}\right)-u\left(q^{\prime}\right)\right|+\left|u\left(q^{\prime}\right)-u(q)\right| .
$$

Then

$$
\left|u(p)-u\left(p^{\prime}\right)\right| \leqslant \int_{\left[p, p^{\prime}\right]}|\nabla u| d x_{1} \quad \text { and } \quad\left|u\left(q^{\prime}\right)-u(q)\right| \leqslant \int_{\left[q^{\prime}, q\right]}|\nabla u| d x_{1}
$$

If $u\left(p^{\prime}\right)=u\left(q^{\prime}\right)$, the assertion follows at once.
If $u\left(p^{\prime}\right) \neq u\left(q^{\prime}\right)$, we can find in [ $\left.p^{\prime}, q^{\prime}\right]$ a shortest interval $\left[p^{\prime \prime}, q^{\prime \prime}\right]$ joining the two closed disjoint subsets $\mathrm{C}_{0} \cap\left[p^{\prime}, q^{\prime}\right]$ and $\mathrm{C}_{1} \cap\left[p^{\prime}, q^{\prime}\right]$. The interior of $\left[p^{\prime \prime}, q^{\prime \prime}\right]$ lies in $D$ and hence

$$
\int_{\left[p^{\prime \prime}, q^{\prime \prime}\right]}|\nabla u| d x_{1}=\mathrm{I}=\left|u\left(p^{\prime}\right)-u\left(q^{\prime}\right)\right| .
$$

This completes the proof of the assertion.
Given now any union E of disjoint intervals $\left[p_{k}, q_{k}\right]$, we have

$$
\begin{aligned}
\sum_{k}\left|u\left(p_{k}\right)-u\left(q_{k}\right)\right| & \leqslant \int_{\mathrm{E} \cap \mathrm{D}}|\nabla u| d x_{1} \leqslant\left(\int_{\mathrm{E}}|\nabla u|^{n} d x_{1}\right)^{\frac{1}{n}}\left(\int_{\mathrm{E}} \mathrm{I} . d x_{1}\right)^{\frac{n-1}{n}} \\
& \leqslant\left(\int_{\mathrm{X} \cap \mathrm{D}}|\nabla u|^{n} d x_{1}\right)^{\frac{1}{n}}(\text { length E })^{\frac{n-1}{n}}
\end{aligned}
$$

Hence $u$ is absolutely continuous on X. It follows that $u$ is ACL everywhere in $\mathbf{R}^{n}$.
It is known that almost all points of a measurable set in $\mathbf{R}^{n}$ are points of linear density in the direction of the coordinate axes (Saks, p. 298). Thus if $\nabla u$ exists at such a point $p$ of $\mathrm{C}_{1} \cup \mathrm{C}_{2}$, we have $\nabla u(p)=0$. Hence $\nabla u=0$ a.e. on $\mathrm{C}_{0} \cup \mathrm{C}_{1}$.

Apart from the smoothly admissible approximations to a given admissible function $u$ on a shell D , there is another kind of approximation that will be needed, namely piecewise linear approximations.
. Let D be an open set such that $\overline{\mathrm{D}} \subset \mathbf{R}^{n}$.
Let $u$ be a continuous function on $\mathbf{R}^{n}$.
Let $\tau$ be a triangulation of $\mathbf{R}^{n}$ into $n$-simplices.
We let $u_{\tau}$ denote the piecewise linear function whose value at a point $\sum_{i=0}^{n} t_{i} p_{i}\left(t_{i}>0, \sum_{i} t_{i}=1\right)$ of a simplex with vertices at $p_{0}, \ldots, p_{n}$ is $\sum_{i=0}^{n} t_{i} u\left(p_{i}\right)$.

We denote by $|\tau|$ the supremum of the diameters of the simplices of $\tau$. Clearly as $|\tau| \rightarrow 0, u_{\tau} \rightarrow u$ uniformly on $\overline{\mathrm{D}}$.

It will be convenient to " standardize " the triangulation $\tau$. Thus for any positive number $a$, we denote by $\tau(a)$ the following triangulation:

Partition $\mathbf{R}^{n}$ into the cubes

$$
\left[k_{1}, \ldots, k_{n}\right]: k_{i} a \leqslant x_{i} \leqslant\left(k_{i}+\mathrm{I}\right) a \quad\left(i=\mathrm{I}, \ldots, n ; k_{i}=\mathrm{o}, \mathrm{I}, \ldots\right),
$$

and then triangulate each cube in the same way.
Finally, there is one other kind of admissible function which will be used. The idea of using such functions goes back to Lebesgue (Rend. Circ. Mat. Palermo, v. 24 (1907), p. 382) who employed them in a paper on Dirichlet's problem of minimizing $\int_{\mathrm{D}}|\nabla u|^{2} d \mathrm{D}$ for plane regions D .

Let D be an open set in a locally connected topological space and let $f$ be a continuous function on $\overline{\mathrm{D}}$. The function $f$ is called monotone on D if for every connected
open subset $\mathrm{U} \subset \mathrm{D}, \sup _{x \in \mathrm{U}} f(x) \leqslant \sup _{y \in \overline{\mathrm{U}}-\mathrm{U}} f(y)$ and $\inf _{x \in \mathrm{U}} f(x) \geqslant \inf _{y \in \overline{\mathrm{U}}-\mathrm{U}} f(y)$ if $\mathrm{U} \neq \overline{\mathrm{U}}$. Thus, for any connected open subset U , if $f$ is constant on the non-empty $\overline{\mathrm{U}}-\mathrm{U}$, then $f$ is constant on U . This property is easily seen to be equivalent to the monotonicity of $f$. And that observation leads to the following method of Lebesgue for " straightening " any function into a monotone function having the same values on $\overline{\mathrm{D}}-\mathrm{D}$.

For any open set U and any point $x$ in U , let $\mathrm{U}_{x}$ denote the connected component of U containing $x$. For any real number $c$, and any continuous function $f$ on D , let $c \mathrm{D} f$ denote the union of all connected components of $\mathrm{D}-f^{-1}(c)$ whose closures lie in D .

For any real number $a$, let $f . a$ denote the function on $\overline{\mathrm{D}}$ defined by

$$
(f . a)(x)=\left\{\begin{array}{ll}
a & \text { if } x \in a \mathrm{D} f \\
f(x) & \text { if } x \notin a \mathrm{D} f
\end{array} .\right.
$$

Thus, $f$ and $f . a$ have the same values on $\overline{\mathrm{D}}--\mathrm{D}$. Moreover, on the boundary of any connected component of $a \mathrm{D} f$, the function $f$ takes on the constant value $a$ and for any non-empty connected subset $\mathrm{T} \ddagger a \mathrm{D} f$, we have $f . a(\mathrm{~T}) \subset f(\mathrm{~T})$ provided $\overline{\mathrm{D}}_{x} \neq \mathrm{D}_{x}$ for all $x \in \mathrm{D}$.

Note. - According to our definition, the only monotone function on a connected topological space X is the constant function. Thus the function $f\left(x_{1}, \ldots, x_{n}\right)=x_{1}$ on $\mathbf{R}^{n}$ is not monotone in our sense. However, it is monotone on any relatively compact open set D. If, in the definition of monotone, we required that $f \mid \mathbf{U}$ take on its maximum and minimum values on $\bar{U}-\mathrm{U}$ for only relatively compact open subsets U , then we would come out with the usual monotone functions on $\mathbf{R}^{1}$. We have not done so since compactness does not enter the proof below. The function $f(x)=x$ on $\mathbf{R}^{1}$ is thus monotone in our sense on all bounded intervals, but not on $\mathbf{R}^{1}$.

Lebesgue's Straightening Lemma. - Let D be an open proper subset of a connected locally connected topological space and let $f$ be a continuous real valued function on $\overline{\mathrm{D}}$ with values in the bounded interval $[m, \mathrm{M}]$. Let $a_{1}, a_{2}, \ldots, a_{n}, \ldots$, be an enumeration of the rational numbers in $[m, \mathrm{M}]$. Set

$$
f_{n}=\left(\ldots\left(\left(f \cdot a_{1}\right) \cdot a_{2}\right) \ldots\right) \cdot a_{n} .
$$

Then $f_{n}$ converges uniformly on $\overline{\mathrm{D}}$ to a monotone function.
Proof. - Let $a$ and $c$ be distinct real numbers and let $x$ be a point of the subset $c \mathrm{D}(f . a)$. Then on the boundary of $(c \mathrm{D}(f . a))_{x}$, the function $f . a$ has the constant value $c$. From the definition of $f . a$ and the fact that $a \neq c$, it follows at once that $f$ has the constant value $c$ on the boundary of $(c \mathrm{D}(f . a))_{x}$ and $x \in c \mathrm{D} f$ if $f(x) \neq c$. Thus $c \mathrm{D}(f . a) \subset c \mathrm{D} f \cup f^{-1}(c)$.

Suppose $a>c$. Then the subset S of elements $y$ in $(c \mathrm{D}(f . a))_{x}$ such that $f(y)>a$ clearly has its boundary in $(c \mathrm{D}(f . a))_{x}$ and hence $\mathrm{S} \subset a \mathrm{D} f \cap c \mathrm{D}(f . a)$. Consequently $(f . a)(y)=a$ for all $y \in \mathrm{~S}$. Equivalently, we can assert

$$
\sup (f . a)(y) \leqslant a
$$

for all $y \in c \mathrm{D}(f . a)$ if $a>c$. Similarly

$$
\inf (f \cdot a)(y) \geqslant a
$$

for all $y \in c \mathrm{D}(f . a)$ if $a<c$.
Clearly $(c \mathrm{D}(f \cdot a))_{x} \neq a \mathrm{D} f$. Hence the values assumed by $f . a$ on $(c \mathrm{D}(f \cdot a))_{x}$ are assumed by $f$ on $(c \mathrm{D}(f \cdot a))_{x}$. Thus

$$
\begin{aligned}
& \sup \{(f . a)(y) ; y \in c \mathrm{D}(f . a)\} \leqslant \sup \{f(y) ; y \in c \mathrm{D} f\} \\
& \inf \{(f . a)(y) ; y \in c \mathrm{D}(f . a)\} \geqslant \inf \{f(y) ; y \in c \mathrm{D} f\} .
\end{aligned}
$$

Given now any distinct real numbers $a_{1}, a_{2}, \ldots, a_{n+1}$, we write

$$
\mathrm{D}_{k}=a_{n+1} \mathrm{D}\left(f \cdot a_{1} \cdot a_{2} \ldots \cdot a_{k}\right)
$$

and we have

$$
\begin{align*}
\sup _{y \in \mathrm{D}_{n}}\left(f \cdot a_{1} a_{2} \ldots a_{n}\right)(y) & \leqslant \sup _{y \in \mathrm{D}_{n-1}}\left(f \cdot a_{1} a_{2} \ldots a_{n-1}\right)(y)  \tag{3}\\
& \leqslant \sup _{y \in \mathrm{D}_{k}}\left(f \cdot a_{1} \cdot a_{2} \ldots a_{k}\right)(y) \\
& \leqslant a_{k} \quad \text { if } a_{k}>a_{n+1} .
\end{align*}
$$

Hence $\sup _{y \in \mathrm{D}_{n}}\left(f_{n}(y)-f_{n+1}(y)\right) \leqslant a_{k}-a_{n+1} \quad$ if $\quad a_{n+1}<a_{k}$.
Similarly,

$$
\inf _{y \in \mathrm{D}_{n}}\left(f_{n}(y)-f_{n+1}(y)\right) \geqslant a_{k}-a_{n+1} \quad \text { if } \quad a_{k}<a_{n+1}
$$

Since $\mathbf{D}_{n}$ is the set on which $f_{n+1}$ differs from $f_{n}$, we find

$$
\sup _{y \in \mathbf{D}}\left|f_{n}(y)-f_{n+1}(y)\right| \leqslant d(n+1) \leqslant d(n)
$$

where $d(n)$ is the length of the longest interval in $[m, \mathrm{M}]-\left\{a_{1}, a_{2}, \ldots, a_{n}\right\}$.
Suppose now that $r$ is an integer less than $n+1$. We shall show by induction on $n-r$ that
(4)

$$
\sup _{y \in \mathbb{D}}\left|f_{r}(y)-f_{n+1}(y)\right| \leqslant d(r) .
$$

Observe first that (4) is true if $n-r=0$, by the assertion above. Next, from $c \mathrm{D}(f . a) \subset c \mathrm{D} f \cup f^{-1}(c)$, we obtain

$$
\begin{aligned}
\mathrm{D}_{n} & \subset \mathrm{D}_{n-1} \cup f_{n}^{-1}\left(a_{n+1}\right) \\
& \subset \mathrm{D}_{r} \cup f_{r+1}^{-1}\left(a_{n+1}\right) \cup \ldots \cup f_{n}^{-1}\left(a_{n+1}\right) \\
& \subset \mathrm{D}_{r} \cup f_{r+1}^{-1}\left(a_{n+1}\right)
\end{aligned}
$$

Since $f_{n+1}$ differs from $f_{n}$ on at most $\mathrm{D}_{n}$, we find by the inductive hypothesis that $\sup _{y \in \mathrm{D}}\left|f_{r}-f_{n+1}\right|$ is the maximum of $d(r)$ and

$$
\sup _{y \in \mathrm{D}_{n}}\left|f_{r}(y)-f_{n+1}(y)\right| .
$$

The latter term is majorized by

$$
\sup _{y \in \mathrm{D}_{r} \cup f_{r}^{1}\left(a_{n+1}\right)}\left|f_{r}(y)-a_{n+1}\right|=\sup _{y \in \mathrm{D}_{r}}\left|f_{r}(y)-a_{n+1}\right| \leqslant d(r)
$$

by the inequality (3) above (and the corresponding inequality for $\inf _{y \in D_{r}} f_{r}$ ). Thus assertion (4) is proved.

Since $\lim _{r \rightarrow \infty} d(r)=0$, it follows immediately that $\left\{f_{n}\right\}$ is uniformly convergent on D to a limit function $\bar{f}$.

The function $\bar{f}$ is monotone on D . Otherwise there would be a real number $c$ such that $c \mathrm{D} \bar{f}$ is not empty. Let $b=\sup _{y} \bar{f}(y)-c$ as $y$ ranges over one of the connected components in $c \mathrm{D} \bar{f}$, say $(c \mathrm{D} \bar{f})_{x}$ and for the sake of definiteness we can suppose that $b=\sup _{y}|\bar{f}(y)-c|, \quad y$ varying over $(c \mathrm{D} \bar{f})_{x} . \quad$ Thus for any $n$ with $\sup \left|f_{n}-\bar{f}\right|<\frac{b}{3}$ we have

$$
\sup \left\{f_{n}(y)-c^{\prime} ; y \in c^{\prime} \mathrm{D} f_{n}\right\}>\frac{b}{3}
$$

whenever $c<c^{\prime}<c+\frac{b}{3}$. Choosing $n$ so that $d(n)>\frac{b}{3}$ and $c<a_{n+1}<c+\frac{b}{3}$, we would have

$$
\sup \left\{f_{n}(y)-a_{n+1} ; y \in a_{n+1} \mathrm{D} f_{n}\right\}>d(n)
$$

that is, $\sup \left|f_{n}-f_{n+1}\right|>d(n)$ contradicting (4).
The proof of the lemma is now complete.

## $\S$ 4. Oscillation of functions with $L^{n}$-integrable gradient

Lemma (4.1). - Let S be the boundary of a ball of radius $r$ in $\mathbf{R}^{n}$ for $n \geqslant 2$, and let $u$ be a function of class $\mathbf{C}^{1}$. Then
where A is a constant depending only on the dimension $n$, but independent of $u$ and $r$, and the gradient is taken along S .

Proof. - Since each side of the inequality is invariant under stretchings, we can assume $r=\mathrm{I}$.

Let $p$ and $q$ be any two distinct points on S .
Let $\pi$ denote the stereographic projection from $p$ onto the tangent plane E to S at the antipodal point of $p$. Set $a=\pi(q), u^{\prime}=u \circ \pi^{-1}, v=|\nabla u| \circ \pi^{-1}$.

Let $\xi$ denote a variable point on S and set $x=\pi(\xi)$. Then $|x|=\tan \theta$, where $\theta$ is defined in the diagram


Since $\pi$ is conformal,

$$
\operatorname{det} \dot{\pi}=\left|\frac{d|x|}{d(2 \theta)}\right|^{n-1}=\left(\frac{\sec ^{2} \theta}{2}\right)^{n-1}=\left(\frac{1+|x|^{2}}{2}\right)^{n-1}
$$

We have

$$
|u(p)-u(q)|=\left|u^{\prime}(\infty)-u^{\prime}(a)\right|=\left|\int_{0}^{\infty} \frac{d u^{\prime}}{d t}(a+t y) d t\right| \leqslant \int_{0}^{\infty}\left|\nabla u^{\prime}\right| d t
$$

where $y$ is any unit vector in E , which we identify with $\mathbf{R}^{n-1}$.
Let Y denote the unit sphere in $\mathbf{R}^{n-1}$, that is the set of all $y$ with $|y|=\mathrm{r}$.
Then

$$
|u(p)-u(q)| \leqslant \frac{\mathrm{I}}{c_{n-2}} \int_{\mathrm{I}} \int_{0}^{\infty}\left|\nabla u^{\prime}\right|(a+t y) d t d y=\frac{\mathrm{I}}{c_{n-2}} \int_{\mathrm{E}}\left|\nabla u^{\prime}\right| \frac{d \mathrm{E}}{|x-a|^{n-2}},
$$

where $c_{n-2}$ is the ( $n-2$ )-measure of the unit $n-2$-sphere. We have seen, for any conformal mapping $\pi$, that

$$
\left|\nabla u^{\prime}\right|=\left(|\nabla u| \circ \pi^{-1}\right)(\operatorname{det} \dot{\pi})^{-\frac{1}{n-1}}
$$

Hence $\left|\nabla u^{\prime}\right|=\frac{2 \eta}{\mathrm{I}+|x|^{2}}$ and we can write

$$
|u(p)-u(q)| \leqslant \frac{2}{c_{n-2}} \int_{\mathrm{E}} \frac{v}{\left(\mathrm{I}+|x|^{2}\right)^{\frac{n-1}{n}}} \cdot \frac{\mathrm{I}}{|x-a|^{n-2}\left(\mathrm{I}+|x|^{2}\right)^{\frac{1}{n}}} d \mathrm{E} .
$$

By Hölder's inequality, the right side is no more than

$$
\frac{2}{c_{n-2}}\left(\int_{\mathrm{E}} \frac{v^{n}}{\left(\mathrm{I}+|x|^{2}\right)^{n-1}} d \mathrm{E}\right)^{\frac{1}{n}}\left(\int_{\mathrm{E}} \frac{\mathrm{I}}{|x-a|^{\frac{(n-2) n}{n-1}}\left(\mathrm{I}+|x|^{2}\right)^{\frac{1}{n-1}}} d \mathrm{E}\right)^{\frac{n-1}{n}}
$$

The denominator of the integral on the right at $x=\infty$ has its order of magnitude $|x|^{-e}$ with

$$
e=\frac{(n-2) n}{n-\mathrm{I}}+\frac{2}{n-\mathrm{I}}=\frac{n^{2}-2 n+2}{n-\mathrm{I}}=\frac{(n-\mathrm{I})^{2}+\mathrm{I}}{n-\mathrm{I}}=n-\mathrm{I}+\frac{\mathrm{I}}{n-\mathrm{I}}>n-\mathrm{I}
$$

and is therefore convergent at $\infty$. At $x=a$, the denominator is essentially $|x-a|^{\frac{(n-2) n}{n-1}}$, the exponent being

$$
\frac{(n-2) n}{n-\mathrm{I}}=\frac{n^{2}-2 n}{n-\mathrm{I}}=\frac{(n-\mathrm{I})^{2}-\mathrm{I}}{n-\mathrm{I}}=n-\mathrm{I}-\frac{\mathrm{I}}{n-\mathrm{I}}<n-\mathrm{I} .
$$

Thus the integral is finite at $x=a$.
Set

$$
\mathrm{A}(a)=\left(\int_{\mathrm{E}}\left[|x-a|^{n-1-\frac{1}{n-1}}\left(\mathrm{I}+|x|^{2}\right)^{\frac{1}{n-1}}\right]^{-1} d \mathrm{E}\right)^{\frac{n-1}{n}}
$$

Then

$$
|u(p)-u(q)|^{n} \leqslant\left(\frac{2 \mathrm{~A}(a)}{c_{n-2}}\right)^{n} \int_{\mathrm{E}} \frac{v^{n}}{\left(\mathrm{I}+|x|^{2}\right)^{n-1}} d \mathrm{E}=\mathrm{A} \int_{\mathrm{S}}|\nabla u|^{n} d \mathrm{~S}
$$

It remains to obtain a bound for $\mathrm{A}(a)$. From $|x| \leqslant|x-a|+|a|$, we see that if $|x-a| \geqslant \frac{1}{2}|a|$, then $|x| \leqslant 3|x-a|$ and thus $|x-a| \geqslant \frac{1}{3}|x|$. If on the other hand $|x-a| \leqslant \frac{|a|}{2}$, then $|x| \geqslant|a|-|x-a| \geqslant \frac{1}{2}|a|$. Now set $b=\sup \left(\frac{1}{2}|a|, \mathrm{I}\right)$. Then write

$$
\int_{\mathrm{E}} \frac{\mathrm{I}}{|x-a|^{n-1-\frac{1}{n-1}}\left(\mathrm{I}+|x|^{2}\right)^{\frac{1}{n-1}}} d \mathrm{E}=\int_{|x-a| \leqslant b}+\int_{|x-a| \geqslant b}
$$

and denote the integrals on the right by I and II.
The first integral, in case $b \leqslant \mathrm{I}$, is bounded by

$$
\begin{aligned}
& \mathrm{I} \leqslant \int_{|x-a| \leqslant b} \frac{\mathrm{I}}{|x-a|^{n-1-\frac{1}{n-1}}} d \mathrm{E}=c_{n-2} \int_{0}^{b} r^{-1+\frac{1}{n-1}} d r=(n-\mathrm{I}) c_{n-2} b^{\frac{1}{n-1}} \\
& \\
& \leqslant(n-\mathrm{I}) c_{n-2} .
\end{aligned}
$$

On the other hand if $b>_{1}$, then $b=\frac{1}{2}|a|$ and $|x| \geqslant b$ on $|x-a| \leqslant b$. Thus

$$
\begin{aligned}
\mathrm{I} & \leqslant\left(\mathrm{I}+b^{2}\right)^{-\frac{1}{n-1}} \int_{|x-a| \leqslant b} \frac{\mathrm{I}}{|x-a|^{n-1-\frac{1}{n-1}}} d \mathrm{E} \leqslant b^{-\frac{2}{n-1}}(n-\mathrm{I}) c_{n-2} b^{\frac{1}{n-1}} \\
& \leqslant(n-\mathrm{I}) c_{n-2} .
\end{aligned}
$$

For the second integral, we have

$$
\begin{aligned}
& I I \leqslant \int_{|x-a| \geqslant b} \frac{1}{\left(\frac{1}{3}|x|\right)^{n-1-\frac{1}{n-1}}\left(I+|x|^{2}\right)^{\frac{1}{n-1}}} d E=\int_{|x-a|>b}^{|x|<1}<\int_{\substack{|x|>1 \\
|x|>b}} \\
& \leqslant 3^{n-1} \int_{|x|<1} \frac{1}{|x|^{n-1-\frac{1}{n-1}}} d \mathrm{E}+3^{n-1} \int_{|x|>1} \frac{1}{|x|^{n-1+\frac{1}{n-1}}} d \mathrm{E} \\
& \leqslant 3^{n-1} c_{n-2} \int_{0}^{1} r^{-1+\frac{1}{n-1}} d r+3^{n-1} c_{n-2} \int_{1}^{\infty} r^{-1-\frac{1}{n-1}} d r=2 \cdot 3^{n-1} c_{n-2}(n-1) .
\end{aligned}
$$

Hence $\quad \mathrm{A} \leqslant \frac{2^{n}}{c_{n-2}^{n}}\left((n-\mathrm{I}) c_{n-2}+2(n-\mathrm{I}) c_{n-2} \cdot 3^{n-1}\right)^{n-1}=\frac{2^{n}(n-1)^{n-1}\left(\mathrm{I}+2 \cdot 3^{n-1}\right)^{n-1}}{c_{n-2}}$.
The proof of Lemma (4.1) is now complete.
As a corollary of Lemma (4.I), we can prove the following result of Loewner on the non-vanishing of conformal capacity.

Lemma (4.2). - Let D be a shell in Moebius space $\mathbf{R}^{n}$ and let $\mathrm{C}_{0}, \mathrm{C}_{1}$ denote the connected components of the complement of D . Then $\mathscr{C}(\mathrm{D})>0$ if neither $\mathrm{C}_{0}$ nor $\mathrm{C}_{1}$ consists of a single point.

Proof. - Choose a point o in $\mathbf{R}^{n}$ such that the two spheres with center at o and radius $r_{1}$ and $r_{2}\left(0<r_{1}<r_{2}\right)$ intersect $\mathrm{C}_{0}$ and $\mathrm{C}_{1}$. Let $\mathrm{S}_{r}$ denote the sphere with center at $o$ and radius $r$. Since $\mathrm{C}_{0}$ and $\mathrm{C}_{1}$ are connected, $\mathrm{S}_{r}$ meets both $\mathrm{C}_{1}$ and $\mathrm{C}_{2}$ for all $r$, $r_{1} \leqslant r \leqslant r_{2}$. Let $\mathrm{D}^{\prime}$ denote the spherical shell: $r_{1} \leqslant|x| \leqslant r_{2}$. Let $u$ be any smoothly admissible function in D . Then

$$
\int_{\mathrm{D}}|\nabla u|^{n} d \mathrm{D}=\int_{\mathbf{R}^{n}}|\nabla u|^{n} d x \geqslant \int_{\mathrm{D}^{\prime}}|\nabla u|^{n} d x=\int_{r_{1}}^{r_{2}} \int_{\mathrm{S}_{r}}|\nabla u|^{n} d \sigma d r
$$

where $d \sigma$ denotes ( $n-1$ )-measure on $\mathrm{S}_{r}$.
Now the gradient $\nabla u$ at any point $x$ in $\mathrm{S}_{r}$ has length no less than the gradient along $\mathrm{S}_{r}$. Hence by Lemma (4. I)

$$
\int_{\mathrm{S}_{r}}|\nabla u|^{n} d \sigma \geqslant \mathrm{~A}^{-1} r^{-1}\left(\operatorname{osc}_{\mathrm{S}_{T}} u\right)^{n}=\mathrm{A}^{-1} r^{-1}
$$

and thus

$$
\int_{\mathrm{D}}|\nabla u|^{n} d \mathrm{D} \geqslant \mathrm{~A}^{-1} \int_{r_{1}}^{r_{2}} \frac{d r}{r}=\mathrm{A}^{-1} \log \frac{r_{2}}{r_{1}}
$$

for all smoothly admissible functions $u$. Taking the infimum over such $u$, we find

$$
\mathscr{C}(\mathrm{D}) \geqslant \mathrm{A}^{-1} \log \frac{r_{2}}{r_{1}}>0
$$

Lemma (4.3). - Let $u$ be continuous and ACL in the shell $\mathrm{D}: a_{1}^{\prime}<\| x_{i}^{\prime} \mid<b$ of $\mathbf{R}^{n}$. Then

$$
\int_{a}^{b}(\underset{|x|=r}{\operatorname{OSC}} u)^{n} \frac{d r}{r} \leqslant \mathrm{~A} \int_{\mathrm{D}}|\nabla u|^{n} d \mathrm{D}
$$

Proof. - We can assume that $|\nabla u|^{n}$ is integrable on D , otherwise the result is trivially true. From Hölder's inequality it follows that $\left|\nabla u_{i}^{\prime}\right|$ is integrable (cf. Proof of Lemma (3.3)). Let $\mathrm{D}^{\prime}$ denote the shell $a+\varepsilon<|x|<b-\varepsilon$ and set

$$
v(x)=\frac{\mathrm{I}}{m(\mathrm{U})} \int_{\mathrm{U}} u(x+y) d y, \quad x \in \mathrm{D}^{\prime}
$$

where $\mathrm{U}=\{y ;|y|<\varepsilon / 2\}$. Then by Lemma 3.3,

$$
\nabla v(x)=\frac{\mathrm{I}}{m(\mathrm{U})} \int_{\mathrm{U}} \nabla u(x+y) d y
$$

and $v$ is of class $\mathrm{C}^{1}$ on $\mathrm{D}^{\prime}$. Hence by Lemma (4.I)

$$
\int_{a+\varepsilon}^{b-\varepsilon}(\underset{|x|=r}{\operatorname{osc}} v)^{n} \frac{d r}{r} \leqslant \mathrm{~A} \int_{D^{\prime}}|\nabla v|^{n} d x \leqslant \mathrm{~A} \int_{D}|\nabla u|^{n} d x
$$

by Lemma (3.3).

Thus

$$
\int_{a+\varepsilon}^{b-\varepsilon}\left(\operatorname{osc}_{|x|=r} v\right)^{n} \frac{d r}{r} \leqslant \mathrm{~A} \int_{\mathrm{D}}|\nabla u|^{n} d \mathrm{D} .
$$

As $\varepsilon \rightarrow 0, v$ converges to $u$ uniformly and hence

$$
\int_{a}^{b}(\underset{|x|=r}{\operatorname{osc}} u)^{n} \frac{d r}{r} \leqslant \mathrm{~A} \int_{\mathrm{D}}|\nabla u|^{n} d \mathrm{D} .
$$

Lemma (4.4). - Let $u$ be continuous and ACL in the half-space $\mathrm{E}_{+}: x_{n}>0$ of $\mathbf{R}^{n}$, and let $\mathrm{S}_{+}$be the hemisphere: $|x|=r, x_{n}>0$. Then

$$
\int_{0}^{\infty}\left(\underset{\mathrm{S}_{+}}{\operatorname{osc} u)^{n}} \frac{d r}{r} \leqslant 2 \mathrm{~A} \int_{\mathrm{E}_{+}}|\nabla u|^{n} d x .\right.
$$

Proof. - For each $\varepsilon>0$, let $\eta$ be the function on $\mathbf{R}^{n}$ :

$$
v\left(x_{1}, x_{2}, \ldots, x_{n}\right)=u\left(x_{1}, x_{2}, \ldots, \varepsilon+\left|x_{n}-\varepsilon\right|\right)
$$

Then $v$ equals $u$ on $x_{n} \geqslant \varepsilon$ and is symmetric with respect to the plane $x_{n}=\varepsilon$.
Thus $v$ is continuous and ACL in $\mathbf{R}^{n}$. By the lemma above,

$$
\int_{0}^{\infty}\left(\operatorname{OSC}_{|x|=r} v\right)^{n} \frac{d r}{r} \leqslant \mathrm{~A} \int_{\mathbf{R}^{n}}|\nabla v|^{n} d x \leqslant 2 \mathrm{~A} \int_{\mathrm{E}_{+}}|\nabla u|^{n} d x .
$$

For each $r>0$, we have $\underset{|x|=r}{\text { osc }} v$ monotone increasing as $\varepsilon$ decreases and

$$
\underset{\mathrm{S}_{+}}{\operatorname{osc}} u=\lim _{\varepsilon \rightarrow 0} \operatorname{osc}_{|x|=r} v .
$$

Hence

$$
\int_{0}^{\infty}\left(\underset{S_{+}}{\operatorname{Osc}} u\right)^{n} \frac{d r}{r}=\lim _{\varepsilon \rightarrow 0} \int_{0}^{\infty}(\underset{|x|=r}{\operatorname{osc}} v)^{n} \frac{d r}{r} \leqslant 2 \mathrm{~A} \int_{\mathrm{E}_{+}}|\nabla u|^{n} d x
$$

 $\mathrm{M}=\int_{\mathrm{D}}|\nabla u|^{n} d \mathrm{D}$. Let $\mathrm{C}_{0}, \mathrm{C}_{1}$ denote the bounded and unbounded components of the complement of D , and extend $u$ to be o on $\mathrm{C}_{0}$ and I on $\mathrm{C}_{1}$. Set

$$
\begin{aligned}
a & =\text { diameter of } \mathrm{C}_{0} \\
b & =d\left(\mathrm{C}_{0}, \mathrm{C}_{1}\right) \\
c(x, y) & =\inf \left(d\left(x, \mathrm{C}_{0}\right), d\left(y, \mathrm{C}_{0}\right)\right) .
\end{aligned}
$$

Then

$$
|u(x)-u(y)|^{n} \leqslant \mathrm{AM}\left(\log \frac{a}{d(x, y)}\right)^{-1}
$$

for all $x, y$ in $\mathbf{R}^{n}$ with $d(x, y)<a$, and

$$
|u(x)-u(y)|^{n} \leqslant \operatorname{AM}\left(\log \frac{c(x, y)}{b}\right)^{-1}
$$

for all $x, y$ with $c(x, y)>b$.

Proof. - By hypothesis, $u$ is monotone on D and therefore on every open subset of $D$. Given any connected open subset $V$ in $\mathbf{R}^{n}$ not contained in or containing $C_{0}$ or $C_{1}$, set $W=V-\left(C_{1} \cup C_{0}\right)$. Clearly the boundary of $W$ meets $C_{i}$ if $V$ does (because $\mathrm{V} \nsubseteq \mathrm{C}_{i}$ and V is connected) and thus

$$
\sup _{x \in V} u(x)=\sup _{x \in \mathrm{~W}} u(x), \quad \inf _{x \in \mathrm{~V}} u(x)=\inf _{x \in \mathrm{~W}} u(x)
$$

because $u$ is constant on each $\mathrm{G}_{i}(i=\mathrm{o}, \mathrm{I})$. Furthermore

$$
u(\overline{\mathrm{~V}}--\mathrm{V})=u(\overline{\mathrm{~W}}-\mathrm{W})
$$

since $\overline{\mathrm{V}}-\mathrm{V} \subset(\overline{\mathrm{W}}-\mathrm{W}) \cup \mathrm{C}_{0} \cup \mathrm{C}_{1}$ and the boundary of V meets $\mathrm{C}_{i}$ if V does (because $\mathrm{C}_{i} \nleftarrow \mathrm{~V}$ and $\mathrm{C}_{i}$ is connected). Thus

$$
\sup _{x \in V} u(x)=\sup _{x \in \overline{\mathrm{~V}}-\mathrm{V}} u(x), \quad \inf _{x \in \mathrm{~V}} u(x)=\inf _{x \in \overline{\mathrm{~V}}-\mathrm{V}} u(x)
$$

because $u$ is monotone on W . It follows that $u$ is monotone on any connected open subset $V$ not containing $\mathrm{C}_{0}$ or $\mathrm{C}_{1}$.

We apply this observation to the open ball V with center at $x$ and radius $a$. The oscillation of $u$ on the sphere $\mathrm{S}_{r}:|z-x|=r$ increases monotonically with $r, o \leqslant r \leqslant a$. We have therefore $|u(x)-u(y)| \leqslant \underset{\mathrm{S}_{r}}{\operatorname{osc} u}$ for $d(x, y)<r<a$ and thus, by Lemma (4.3)

$$
|u(x)-u(y)|^{n} \int_{d(x, y)}^{a} \frac{d r}{r} \leqslant \mathrm{AM}
$$

On the other hand, we can apply the observation above to the exterior V of the closed ball with center $p$ and radius $r>d\left(p, \mathrm{C}_{1}\right)$ where $p$ is a point in $\mathrm{C}_{0}$ with $d\left(p, \mathrm{C}_{1}\right)=d\left(\mathrm{C}_{0}, \mathrm{C}_{1}\right)=b$. The fact that $u$ is monotone on V implies here that
 $|u(x)-u(y)| \leqslant \underset{S_{r}}{\operatorname{osc}} u$ if $r<d(p, x), d(p, y)$. By Lemma (4.3)

$$
|u(x)-u(y)|^{n} \int_{b}^{c(x, y)} \frac{d r}{r} \leqslant \mathrm{AM}
$$

These are the inequalities of our lemma.

## § 5. EXTREMAL ADMISSIBLE FUNCTIONS

The following lemma is due to J. A. Clarkson (On uniformly convex spaces, Trans. Amer. Math. Soc., v. $4^{\circ}$ (1936), pp. 396-414).

Lemma (5.1). - For any two $\mathbf{C}^{n}$-valued functions $x$ and $y$ on a measure space E whose component functions are in $\mathrm{L}^{p}, p \geqslant 2$,

$$
\begin{equation*}
2\left(\|x\|^{p}+\|y\|^{p}\right) \leqslant\|x+y\|^{p}+\|x-y\|^{p} \leqslant 2^{p-1}\left(\|x\|^{p}+\|y\|^{p}\right) \tag{I}
\end{equation*}
$$

where $\|x\|$ denotes $\||x|\|_{p}=\left(\int_{E}\left(\left|x_{1}\right|^{2}+\ldots+\left|x_{n}\right|^{2}\right)^{p / 2}\right)^{\frac{1}{p}}$.
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Proof. - Observe first that the right inequality in (1) is equivalent to the left inequality; to see this, set $x^{\prime}=x+y, y^{\prime}=x-y$. Thus it suffices to prove only the left inequality.

First we establish the left inequality when E is a point; that is

$$
2\left(|x|^{p}+|y|^{p}\right) \leqslant|x+y|^{p}+|x-y|^{p}
$$

for any two complex $n$-tuples $x$ and $y$. We may assume without loss of generality that $|y| \leqslant|x|$. Dividing by $|x|$, the inequality ( $\mathrm{I}^{\prime}$ ) is equivalent to

$$
2\left(\mathrm{I}+|c|^{p}\right) \leqslant|u+c|^{p}+|u-c|^{p}
$$

for any complex $n$-tuples $u$ and $c$ with $|u|=\mathrm{I},|c| \leqslant \mathrm{I}$. Set $2 a=|u+c|+|u-c|$, $\mathrm{I}+b=\sup \left(\frac{|u+c|}{a}, \frac{|u-c|}{a}\right) . \quad$ Then $2 a \geqslant|2 u|=2, b \geqslant 0$, and $\mathrm{I}-b=\inf \left(\frac{|\mathrm{I}+c|}{a}, \frac{|\mathrm{I}-c|}{a}\right)$. Hence $|u+c|^{p}+|u-c|^{p}=a^{p}\left((\mathrm{I}+b)^{p}+(\mathrm{I}-b)^{p}\right)$ where $a \geqslant \mathrm{I}$ and $0 \leqslant b \leqslant \mathrm{I}$. For $p=2$, we have equality in ( $\mathrm{I}^{\prime}$ ) and ( $\left.\mathrm{I}^{\prime \prime}\right)$. As $p$ increases, the left side of ( $\mathrm{I}^{\prime \prime}$ ) decreases or stays constant. On the other hand, the right side of ( $r^{\prime \prime}$ ) increases with $p$, for the derivative of $f(t)=(\mathrm{I}+b)^{t}+(\mathrm{I}-b)^{t}$ is, for $t \geqslant \mathrm{I}$,

$$
\begin{aligned}
f^{\prime}(t) & =(\mathrm{I}+b)^{t} \log (\mathrm{I}+b)+(\mathrm{I}--b)^{t} \log (\mathrm{I}-b) \geqslant(\mathrm{I}+b) \log (\mathrm{I}+b)+(\mathrm{I}-b) \log (\mathrm{I}-b) \\
& =-2\left(\frac{b^{2}}{2}+\frac{b^{4}}{4}+\ldots\right)+2 b\left(b+\frac{b^{3}}{3}+\ldots\right) \\
& =2\left(b^{2}\left(\mathrm{I}-\frac{\mathrm{I}}{2}\right)+b^{4}\left(\frac{\mathrm{I}}{3}-\frac{\mathrm{I}}{4}\right)+\ldots\right)>0 .
\end{aligned}
$$

Hence ( $\mathrm{I}^{\prime \prime}$ ) and therefore ( $\mathrm{I}^{\prime}$ ) is established.
Integrating the inequality ( $I^{\prime}$ ) which holds at each point for any $\mathbf{C}^{n}$-valued functions on $E$, we obtain the left side of (I).

Theorem (5.1). - Let D be a shell whose boundary components each contains more than one point. Then there exists a unique admissible function u for D such that

$$
\mathscr{C}(\mathrm{D})=\int_{\mathrm{D}}|\nabla u|^{n} d \mathrm{D}
$$

The function $u$ is monotone on D .
Proof. - Let $\mathrm{C}_{0}$ and $\mathrm{C}_{1}$ denote the components of the complement of D in Moebius space. Since $\mathscr{C}(\mathrm{D})$ is invariant under conformal mappings, we can assume without loss of generality that $\infty \in \mathrm{C}_{1}$. The hypothesis of the theorem assures that the constants $a=$ diameter $\mathrm{C}_{0}$ and $b=d\left(\mathrm{C}_{0}, \mathrm{C}_{1}\right)$ of Lemma (4.5) are both positive. The inequalities of Lemma (4.5) then imply that the family of monotone admissible functions $u$ such that $\int_{\mathrm{D}}|\nabla u|^{n} d \mathrm{D} \leqslant \mathrm{M}$ is an equicontinuous family of functions on the entire Moebius space $\mathbf{R}^{n} \cup \infty$ (without further warning we shall denote by the same letter an admissible function and its trivial extension to $\mathbf{C}_{\mathbf{0}}$ and $\mathbf{C}_{\mathbf{1}}$ ). Given therefore any sequence of func-
tions in this family, we can, by Ascoli's theorem, select a uniformly convergent subsequence.

Consider now any sequence $u_{1}, u_{2}, \ldots$, of monotone admissible functions such that

$$
\lim _{r \rightarrow \infty} \int_{\mathrm{D}}\left|\nabla u_{r}\right|^{n} d \mathrm{D}=\mathscr{C}(\mathrm{D})
$$

By Clarkson's Lemma, we have

$$
\begin{equation*}
\left\|\frac{\nabla u_{r}+\nabla u_{s}}{2}\right\|^{n}+\left\|\frac{\nabla u_{r}-\nabla u_{s}}{2}\right\|^{n} \leqslant \frac{1}{2}\left(\left\|\nabla u_{r}\right\|^{n}+\left\|\nabla u_{s}\right\|^{n}\right) . \tag{2}
\end{equation*}
$$

Since $\frac{\mathrm{I}}{2}\left(u_{r}+u_{s}\right)$ is an admissible function for D , it follows from (2) that

$$
\lim _{r, s \rightarrow \infty} \int_{\mathrm{D}}\left|\nabla\left(\frac{u_{r}+u_{s}}{2}\right)\right|^{n} d \mathrm{D}=\mathscr{C}(\mathrm{D})
$$

and that

$$
\lim _{r, s \rightarrow \infty}\left\|\nabla u_{r}-\nabla u_{s}\right\|=0
$$

Hence the $i$-th component function of $\nabla u_{r}$ converges in $\mathrm{L}^{n}$ to a function $f_{i}(i=\mathrm{I}, \ldots, n)$. We may suppose, by the remark above, that in fact the sequence $u_{1}, u_{2}, \ldots$ converges uniformly on Moebius space to a limit function $u$. We show that $u$ is admissible for D , $\nabla u=\left(f_{1}, \ldots, f_{n}\right)$ a.e., and $\int_{\mathrm{D}}|\nabla u|^{n} d \mathrm{D}=\mathscr{C}(\mathrm{D})$.

Let $\mathrm{I}=[p, q]$ be a finite line segment in $\mathbf{R}^{n}$ parallel to the $x_{1}$-axis. Let $\mathrm{J}_{\varepsilon}$ denote the cube $\left|x_{2}\right|<\varepsilon, \ldots\left|x_{n}\right|<\varepsilon$ in the $(n-\mathrm{I})$-plane $x_{1}=0$. Since $\frac{\partial u_{r}}{\partial x_{1}}$ converges in $\mathrm{L}^{n}$ to $f_{1}$, it converges, by Hölder's inequality, in $L^{1}$ to $f_{1}$ on any bounded region and thus

$$
\begin{aligned}
\int_{I \times J_{\varepsilon}} f_{1}(x+y) d x_{1} d y & =\lim _{r \rightarrow \infty} \int_{\mathrm{I} \times \mathrm{J}_{\varepsilon}} \frac{\partial u_{r}}{\partial x_{1}} d x_{1} d y \\
& =\lim _{r \rightarrow \infty}\left(\int_{J_{\Sigma}} u_{r}(q+y) d y-\int_{J_{\varepsilon}} u_{r}(p+y) d y\right) \\
& =\int_{J_{\varepsilon}}(u(q+y)-u(p+y)) d y .
\end{aligned}
$$

Dividing by $\varepsilon^{n-1}$ and letting $\varepsilon \rightarrow 0$, the left side approaches $\int_{\mathrm{I}} f_{1}(x) d x_{1}$ for almost all $\left(x_{2}, \ldots, x_{n}\right)$, whereas the right side approaches $u(q)-u(p)$, since $u$ is continuous. Thus $\int_{p}^{q} f_{1} d x_{1}=u(q)-u(p), u$ is absolutely continuous on almost all line segments [ $p, q$ ] parallel to the $x_{1}$-axis, and $\frac{\partial u}{\partial x_{1}}=f_{1}$ a.e. This argument applies to lines parallel to any of the coordinate axes and thus $u$ is ACL and $\nabla u=f=\left(f_{1}, \ldots, f_{n}\right)$ a.e. Consequently,

$$
\int_{\mathrm{D}}|\nabla u|^{n} d \mathrm{D}=\|f\|^{n}=\lim _{r \rightarrow \infty}\left\|\nabla u_{r}\right\|^{n}=\mathscr{C}(\mathrm{D})
$$

Given any admissible function $v$ on D such that $\int_{\mathbf{D}}|\nabla u|^{n} d \mathrm{D}=\mathscr{C}(\mathrm{D})$, we obtain by Clarkson's inequality

$$
\left\|\left.\nabla\left(\frac{u+v}{2}\right)\right|^{n}+\right\| \frac{\nabla u-\nabla v}{2} \|^{n} \leqslant \mathscr{C}(\mathrm{D})
$$

and, since $\frac{u+v}{2}$ is admissible for $\mathrm{D}, \nabla(u-v)=0$ a.e. It follows, since $u$ and $v$ are continuous, that $u-v=0$. This establishes the uniqueness of $u$.

The fact that $u$ is monotone follows at once from the observation that Lebesgue's straightening process would yield an admissible function $u_{0}$ with $\int_{D}\left|\nabla u_{0}\right|^{n} d \mathrm{D}<\int_{D}|\nabla u|^{n} d \mathrm{~S}$ if $u$ were not monotone.

Definition. - The function u of Theorem 5.I is called the extremal function for D. (It depends of course on which of the components of the complement of D is labelled $\mathrm{C}_{0}$.)

Note. - The extremal function satisfies the variational condition

$$
\int_{\mathrm{D}}|\nabla u|^{n-2} \nabla u \cdot \nabla w d \mathrm{D}=\mathrm{o}
$$

for any $\mathrm{C}^{1}$ function $w$ having compact support in D , or equivalently, in the weak sense

$$
\int_{\mathrm{D}}\left(\operatorname{div}|\nabla u|^{n-2} \nabla u\right) w d \mathrm{D}=0 .
$$

Thus $\operatorname{div}|\nabla u|^{n-2} \nabla u=0$ in the weak sense.
Using in succession results of E. di Giorgi, C. B. Morrey, and E. Hopf on elliptic partial differential equations, F. W. Gehring proves for $n=3$ in "Rings and quasiconformal mappings in space ", Trans. Amer. Math. Soc., v. 103 (196), that if $\frac{\mathrm{I}}{\mathrm{M}}<\nabla u<\mathrm{M}$ a.e., then $u$ is real analytic. Gehring's proof applies with obvious modifications for $n \geqslant 2$ as well.

We shall not require the analyticity theorem in this account.

## § 6. CONTINUITY OF $\mathscr{C}(\mathrm{D})$ IN D

We introduce a metric topology on the closed subsets of Moebius space $\mathbf{R}^{n} \cup \infty$, defining the distance between two closed subsets $\mathbf{E}$ and F :

$$
p(\mathrm{E}, \mathrm{~F})=\sup _{x \in \mathrm{E}-\mathrm{F}} p(x, \mathrm{~F})+\sup _{x \in \mathrm{~F}-\mathrm{E}} p(x, \mathrm{E})
$$

$p(x, y)$ denoting the standard metric on the sphere $\mathrm{S}^{n}$ which is identified with Moebius space via some fixed stereographic projection. This metric topology is equivalent to the topology obtained from the standard euclidean distance $d(x, y)$ when we restrict attention to the compact subsets of $\mathbf{R}^{n}$. Without this restriction, the latter topology is finer.

For arbitrary shells D and $\mathrm{D}^{\prime}$, we define the distance between D and $\mathrm{D}^{\prime}$
as $p\left(\mathrm{C}_{0} \cup \mathrm{C}_{1}, \mathrm{C}_{0}^{\prime} \cup \mathrm{C}_{1}^{\prime}\right)$ where $\mathrm{C}_{0} \cup \mathrm{C}_{1}$ and $\mathrm{C}_{0}^{\prime} \cup \mathrm{C}_{1}^{\prime}$ are the complements of D and $\mathrm{D}^{\prime}$ respectively.

Theorem (6.1). - $\mathscr{C}(\mathrm{D})$ depends continuously on D .
Proof. - We must show that given any $\varepsilon>0$, we have $\left|\mathscr{C}(\mathrm{D})-\mathscr{C}\left(\mathrm{D}^{\prime}\right)\right|<\varepsilon$ for all shells $\mathrm{D}^{\prime}$ in some neighborhood of D . Let $x_{1}$ denote a fixed point in $\mathrm{C}_{1}$, and let I denote the set of all shells not containing $x_{1}$. For any neighborhood $U$ of $D$ in the space of all shells, the orbit under the Moebius group of $\mathrm{U} \cap \mathrm{I}$ is a neighborhood of D . For, since the metric as well as conformal capacity is preserved by a transitive subgroup of the Moebius group, we can assume $\infty \in \mathrm{D}$. Let $\mathrm{U}_{\delta}$ denote the set of all shells $\mathrm{D}^{\prime}$ such that $d\left(\mathrm{C}_{0}^{\prime} \cup \mathrm{C}_{1}^{\prime}, \mathrm{C}_{0} \cup \mathrm{C}_{1}\right)<\delta$. Given $\mathrm{D}^{\prime} \in \mathrm{U}_{\delta}$, we can by a translation $t$ of length at most $\delta$ carry $\mathrm{C}_{1}^{\prime}$ into a subset containing the point $x_{1}$ of $\mathrm{G}_{1}$. Since

$$
d\left(\mathrm{C}_{0} \cup \mathrm{C}_{1}, t\left(\mathrm{C}_{0}^{\prime} \cup \mathrm{C}_{1}^{\prime}\right)\right) \leqslant d\left(\mathrm{C}_{0} \cup \mathrm{C}_{1}, \mathrm{C}_{0}^{\prime} \cup \mathrm{C}_{1}^{\prime}\right)+d\left(\mathrm{C}_{0}^{\prime} \cup \mathrm{C}_{1}^{\prime}, t\left(\mathrm{C}_{0}^{\prime} \cup \mathrm{C}_{1}^{\prime}\right)\right)<3 \delta
$$

we have $\mathrm{T}\left(\mathrm{U}_{3 \delta} \cap \mathrm{I}\right) \supset \mathrm{U}_{\delta}$, where T denotes the group of translations on $\mathbf{R}^{n} \cup \infty$. Since $\mathrm{C}_{0} \cup \mathrm{C}_{1} \subset \mathbf{R}^{n}$, the subsets $\mathrm{U}_{\delta}$ form a base of neighborhoods of D , and thus our assertion is established.

Since $\mathscr{C}(\mathrm{D})$ is invariant under the Moebius group, it suffices to show that $\mathscr{C}$ is continuous on I. We may clearly take $x_{1}$ to be any point in Moebius space. We now fix $x_{1}=\infty$. That is, it remains only to prove that $\mathscr{C}(D)$ depends continuously on $D$ for shells $\mathrm{D} \subset \mathbf{R}^{n}$.

Let D be a shell in $\mathbf{R}^{n}$, let $\varepsilon>0$, and let $u$ be a smoothly admissible function for D such that $\int_{\mathrm{D}}|\nabla u|^{n} d \mathrm{D}<\mathscr{C}(\mathrm{D})+\varepsilon$. Let K denote the support of $\nabla u$. By definition of $u, \mathrm{~K}$ is a compact subset of D . For any shell $\mathrm{D}^{\prime}$ sufficiently near D , we have $\mathrm{K} \subset \mathrm{D}^{\prime}$. Hence $u$ is admissible for $D^{\prime}$ and

$$
\begin{equation*}
\mathscr{C}\left(\mathrm{D}^{\prime}\right)<\mathscr{C}(\mathrm{D})+\varepsilon \tag{I}
\end{equation*}
$$

The opposite inequality will be deduced effectively from Lemma (4.5) which provides uniform estimates for the variation of monotone admissible functions. Let $\mathrm{C}_{0}$, $\mathrm{C}_{0}^{\prime}$ denote the bounded component of the complement of D and $\mathrm{D}^{\prime}$ respectively. Let $a^{\prime}$ denote the diameter of $\mathrm{C}_{0}^{\prime}$ in the euclidean metric, let $b^{\prime}=d\left(\mathrm{C}_{0}^{\prime}, \mathrm{C}_{1}^{\prime}\right)$, and $c^{\prime}(x, y)=\inf \left(d\left(x, \mathrm{C}_{0}^{\prime}\right), d\left(y, \mathrm{C}_{0}^{\prime}\right)\right)$. For each shell $\mathrm{D}^{\prime}$ in $\mathbf{R}^{n}$, select a monotone admissible function $u^{\prime}$ such that

$$
\int_{\mathrm{D}^{\prime}}\left|\nabla u^{\prime}\right|^{n} d \mathrm{D}^{\prime} \leqslant \mathscr{C}\left(\mathrm{D}^{\prime}\right)+\varepsilon
$$

(We could for example take $u^{\prime}$ to be the extremal function.) Define $u^{\prime}$ to be constant on $\mathrm{C}_{0}^{\prime}$ and $\mathrm{C}_{1}^{\prime}$.

Define the function $v^{\prime}$ on $\mathbf{R}^{n}$ :

$$
v^{\prime}(x)= \begin{cases}0 & \text { if } u^{\prime}(x)<\varepsilon \\ \frac{u^{\prime}(x)-\varepsilon}{\mathrm{I}-2 \varepsilon} & \text { if } \quad \varepsilon \leqslant u^{\prime}(x) \leqslant \mathrm{I}-\varepsilon \\ \mathrm{I} & \text { if } \mathrm{I}-\varepsilon<u^{\prime}(x) \leqslant \mathrm{I}\end{cases}
$$

Then $v^{\prime}$ is ACL in $\mathbf{R}^{n}$. Let $\mathrm{K}^{\prime}$ denote the support of $\nabla v^{\prime}$.
Then $\mathrm{K}^{\prime}$ is a compact subset of $\mathrm{D}^{\prime}$. The inequality

$$
\left|u^{\prime}(x)-u^{\prime}(y)\right|^{n} \leqslant \operatorname{AM}\left(\log \frac{a^{\prime}}{d(x, y)}\right)^{-1} \quad \text { if } d(x, y)<a^{\prime}
$$

implies for any $x \in \mathrm{~K}^{\prime}$,

$$
d\left(x, \mathrm{C}_{i}^{\prime}\right) \geqslant \frac{a^{\prime}}{\exp \left(\mathrm{AM}^{\prime} \varepsilon^{-n}\right)} \geqslant \frac{a-\varepsilon}{\exp \left(\mathrm{A}(\mathrm{M}+2 \varepsilon) \varepsilon^{-n}\right)}=\delta(\varepsilon)
$$

( $i=0, \mathrm{I}$ ), provided $\mathrm{D}^{\prime}$ is sufficiently close to D .
The inequality

$$
\left|u^{\prime}(x)-u^{\prime}(y)\right|^{n} \leqslant \mathrm{AM}^{\prime}\left(\log \frac{c^{\prime}(x, y)}{b}\right)^{-1} \quad \text { if } \quad c^{\prime}(x, y)<a^{\prime}
$$

implies for any $x \in \mathrm{~K}^{\prime}$ (when we take $y=\infty$ )

$$
d\left(x, \mathrm{C}_{0}^{\prime}\right) \leqslant b^{\prime} \exp \left(\mathrm{AM}^{\prime} \varepsilon^{-n}\right)<(b+\varepsilon) \exp \left(\mathrm{A}(\mathbf{M}+2 \varepsilon) \varepsilon^{-n}\right)=r(\varepsilon)
$$

provided $\mathrm{D}^{\prime}$ is sufficiently close to D .
Thus for any shell $\mathrm{D}^{\prime}$ satisfying

$$
\begin{equation*}
d\left(\left(\mathrm{C}_{0}^{\prime} \cup \mathrm{C}_{1}^{\prime}\right) \cap \mathrm{B},\left(\mathrm{C}_{0} \cup \mathrm{C}_{1}\right) \cap \mathrm{B}\right)<\delta(\varepsilon) \tag{2}
\end{equation*}
$$

where $B$ is the closed $n$-ball with center at some fixed point of $\mathrm{C}_{0}$ and radius $r(\varepsilon)+b(\varepsilon)+\delta(\varepsilon)$, and for any $x \in \mathrm{~K}^{\prime}$, we have $d\left(x, \mathrm{C}_{0} \cup \mathrm{C}_{1}\right)>0$; that is, $\mathrm{K}^{\prime} \subset \mathrm{D}$.

Consequently, $v^{\prime}$ is admissible for D and therefore

$$
\begin{equation*}
\mathscr{C}(\mathrm{D}) \leqslant \int_{\mathrm{D}}(\nabla v)^{n} d \mathrm{D} \leqslant(\mathrm{I}-2 \varepsilon)^{-n}\left(\mathscr{C}\left(\mathrm{D}^{\prime}\right)+\varepsilon\right) \tag{3}
\end{equation*}
$$

for all shells $\mathrm{D}^{\prime}$ satisfying the condition (2). The set of such shells is a neighborhood of D. From assertions (1) and (3) the theorem follows.

Remark. - For closed subsets of $\mathbf{R}^{n} \cup \infty$ containing $\infty$, the metric topology introduced above is equivalent to the inductive topology with respect to compact sets of $\mathbf{R}^{n}$. The same is true therefore for closed subsets of $\mathbf{R}^{n} \cup \infty$. Let D and $\mathrm{D}_{m}$ denote shells in $\mathbf{R}^{n} \cup \infty$ with complementary components $\mathrm{C}_{0}, \mathrm{C}_{1}$ and $\mathrm{C}_{m, 0}, \mathrm{C}_{m, 1}$ respectively such that

$$
\begin{aligned}
& \lim _{m \rightarrow \infty} \mathrm{C}_{m, 0} \cap \mathrm{~K}=\mathrm{C}_{0} \cap \mathrm{~K} \\
& \lim _{m \rightarrow \infty} \mathrm{C}_{m, 1} \cap \mathrm{~K}=\mathrm{C}_{1} \cap \mathrm{~K}
\end{aligned}
$$

for each compact set $K \subset \mathbf{R}^{n}$. Then $D_{m}$ converges to D in the topology of our space of shells.

## § 7. THE CONFORMAL CAPACITY AND MODULUS OF SOME SPECIAL SHELLS

Let $\mathrm{D}_{a, b}$ denote the spherical shell consisting of all $x$ in $\mathbf{R}^{n}$ with $a<|x|<b$, where $0<a<b$.

Lemma (7.1). $-\mathscr{C}\left(\mathrm{D}_{a, b}\right)=c_{n-1}\left(\log \frac{b}{a}\right)^{-(n-1)}$, where $c_{n-1}$ denotes the $(n-1)$ measure of the surface S of the unit n-ball.

Proof. - Let $u$ be an admissible $\mathrm{C}^{1}$ function for $\mathrm{D}=\mathrm{D}_{a, b}$. Then integrating along a ray, we get

$$
\begin{aligned}
\mathrm{I} & \leqslant\left(\int_{a}^{b}|\nabla u| d r\right)^{n}=\left(\int_{a}^{b}|\nabla u| r^{\frac{n-1}{n}} r^{-\frac{n-1}{n}} d r\right)^{n} \\
& \leqslant\left(\int_{a}^{b}|\nabla u|^{n} r^{n-1} d r\right)\left(\int_{a}^{b} r^{-1} d r\right)^{n-1}
\end{aligned}
$$

by Hölder's inequality. Integrating over all rays, we get

$$
c_{n-1} \leqslant \int_{\mathrm{D}}|\nabla u|^{n} d \mathrm{D} \cdot\left(\log \frac{b}{a}\right)^{n-1}
$$

Hence $\mathscr{C}(\mathrm{D}) \geqslant c_{n-1}\left(\log \frac{b}{a}\right)^{-(n-1)}$. On the other hand, taking $u(x)=\left(\log \frac{|x|}{a}\right)\left(\log \frac{b}{a}\right)^{-1}$,
we find

$$
\int_{\mathrm{D}}|\nabla u|^{n} d \mathrm{D}=c_{n-1}\left(\log \frac{b}{a}\right)^{-n} \int_{a}^{b}\left(\frac{\mathrm{I}}{r}\right)^{n} r^{n-1} d r=c_{n-1}\left(\log \frac{b}{a}\right)^{-(n-1)}
$$

Hence $\mathscr{C}(\mathbf{D})=c_{n-1}\left(\log \frac{b}{a}\right)^{-(n-1)}$.
It is convenient to introduce another conformal invariant, equivalent to conformal capacity, but rendering somewhat simpler the formulas to be encountered below.

Definition. - The modulus of a shell D in Moebius $n$-space is $\left(\frac{c_{n-1}}{\mathscr{C}(\mathrm{D})}\right)^{\frac{1}{n-1}}$. It is
$d \bmod$. denoted $\bmod \mathrm{D}$.

Thus we have $\bmod \mathrm{D}_{a, b}=\log \frac{b}{a}$, generalizing the familiar formula for $n=2$.
Given two shells $D$ and $D^{\prime}$ such that $C_{0}^{\prime} \supset C_{0}$ and $C_{1}^{\prime} \supset C_{1}$, where $C_{0}$ and $C_{1}$ are two components of the complement of $D$ and $C_{0}^{\prime}$ and $C_{1}^{\prime}$ are the components of the complement of $D^{\prime}$, we say that $D^{\prime}$ separates the boundaries of $D$. It is clear that $D^{\prime} \subset D$ and that any admissible function for $\mathrm{D}^{\prime}$ is admissible for D . Hence

$$
\mathscr{C}(\mathrm{D}) \leqslant \mathscr{C}\left(\mathrm{D}^{\prime}\right) \quad \text { and } \quad \bmod \mathrm{D} \geqslant \bmod \mathrm{D}^{\prime}
$$

if $\mathrm{D}^{\prime}$ separates the boundaries of D .
It follows from this and Lemma (7.1) that $\mathscr{C}(\mathrm{D})=0$ if one of the components of the complement of D reduces to a single point. For without loss of generality, one can assume that the origin is the unique point of the component $\mathrm{C}_{0}$. Then for suitably small posi-
tive $b$, the spherical shell $\mathrm{D}_{a, b}$ is contained in D and it separates the boundaries of D . Hence

$$
\bmod \mathrm{D} \geqslant \log \frac{b}{a}
$$

for arbitrarily small $a$. It follows that $\bmod \mathrm{D}=\infty$ and $\mathscr{C}(\mathrm{D})=0$.
Note. - The condition that the component $\mathrm{C}_{0}$ of the complement of the shell D consist of a single point is equivalent to the condition that the boundary component $\mathrm{C}_{0} \cap \overline{\mathrm{D}}$ consist of a single point $p$. For the complement of a single point in $S^{n}$ is connected. On the other hand $\mathrm{G} \cup \mathrm{D}=\mathrm{S}^{n}-\mathrm{C}_{0}$ is open, and $\mathrm{C}_{0}-p=\mathrm{S}^{n}-\left(\mathrm{C}_{1} \cup \mathrm{D}\right)-p=\mathrm{S}^{n}-\left(\mathrm{C}_{1} \cup \overline{\mathrm{D}}\right) \quad$ is open. Since $\quad \mathrm{S}^{n}-p=\left(\mathrm{C}_{1} \cup \mathrm{D}\right) \cup\left(\mathrm{C}_{0}-p\right)$, we conclude that $\mathrm{C}_{0}-p$ is empty.

Referring to Lemma (4.2), we can now assert that the conformal capacity of a shell is zero if and only if one of its boundary components reduces to a point.

Lemma (7.2). - Let D be a shell and let $\mathrm{D}_{1}, \ldots, \mathrm{D}_{m}$ be mutually disjoint shells each separating the boundaries of D . Then

$$
\bmod \mathrm{D} \geqslant \bmod \mathrm{D}_{1}+\ldots+\bmod \mathrm{D}_{m}
$$

Proof. - By hypothesis we have $\mathrm{C}_{0} \subset \mathrm{C}_{0, i}$ and $\mathrm{C}_{1} \subset \mathrm{C}_{1, i}$ for the connected components of the complements. Let $u_{i}$ be a smoothly admissible function for $\mathrm{D}_{i}$ taking the value $o$ on $\mathrm{C}_{0, i}$ and I on $\mathrm{C}_{1, i}(i=\mathrm{I}, \ldots, m)$. Set $u=\sum_{i=1}^{m} a_{i} u_{i}$, where $a_{i} \geqslant 0$ and $\sum_{i=1}^{m} a_{i}=\mathrm{I}$. Then

$$
\int_{\mathrm{D}}|\nabla u|^{n} d x=\sum_{i=1}^{m} a_{i}^{n} \int_{D_{i}}\left|\nabla u_{i}\right|^{n} d x
$$

Hence $\mathscr{C}(\mathrm{D}) \leqslant \sum_{i=1}^{m} a_{i}^{n} \mathscr{C}\left(\mathrm{D}_{i}\right)$. We can assume $\mathscr{C}(\mathrm{D})>0$, otherwise there is nothing to prove. Hence $\mathscr{C}\left(\mathrm{D}_{i}\right)>0$ and we may take $a_{i}=\mathscr{C}\left(\mathrm{D}_{i}\right)^{-\frac{1}{n-1}}\left(\sum_{j=1}^{m} \mathscr{C}\left(\mathrm{D}_{j}\right)^{-\frac{1}{n-1}}\right)^{-1}$. Then we get

$$
\begin{aligned}
& \mathscr{C}(\mathrm{D}) \leqslant \frac{\sum_{i=1}^{m} \mathscr{C}\left(\mathrm{D}_{i}\right)^{-\frac{1}{n-1}}}{\left(\sum_{i=1}^{m} \mathscr{C}\left(\mathrm{D}_{i}\right)^{-\frac{1}{n-1}}\right)^{n}}=\left(\sum_{i=1}^{m} \mathscr{C}\left(\mathrm{D}_{i}\right)^{\frac{1}{n-1}}\right)^{-(n-1)} \\
& \mathscr{C}(\mathrm{D})^{-\frac{1}{n-1}} \geqslant \sum_{i=1}^{m} \mathscr{C}\left(\mathrm{D}_{i}\right)^{-\frac{1}{n-1}} \\
& \bmod \mathrm{D} \geqslant \sum_{i=1}^{m} \bmod \mathrm{D}_{i}
\end{aligned}
$$

Definition. - The Grötzsch shell $\mathrm{D}_{\mathrm{G}} \mathrm{D}_{\mathrm{G}}=(a)$ is the shell in $\mathbf{R}^{n}$ whose complementary components consist of the sphere $|x|<1$ and the ray $a \leqslant x_{1}<\infty, x_{2}=x_{3}=\ldots=x_{n}=0$, where $a>1$.

The Teichmüller shell $\mathrm{D}_{\mathrm{T}}=\mathrm{D}_{\mathrm{T}}(b)$ is the shell in $\mathbf{R}^{n}$ whose complementary components consist of the segment $-\mathrm{I} \leqslant x_{1} \leqslant 0, x_{2}=\ldots=x_{n}=0$ and the ray $b<x_{1}<\infty$, $x_{2}=x_{3}=\ldots=x_{n}=0$, where $b>0$.

Teichmüller has introduced the functions $\Phi_{n}$ and $\Psi_{n}$ (in the case $n=2$ ) defined by:

$$
\bmod \mathrm{D}_{\mathrm{G}}=\log \Phi_{n}(a), \quad \bmod \mathrm{D}_{\mathrm{T}}=\log \Psi_{n}(b)
$$

Lemma (7.3). - $\frac{\Phi_{n}(a)}{a}$ is a monotonically increasing function of $a, a>1$. For any $b>0$,

$$
\Psi_{n}(b)=\left(\Phi_{n}((b+\mathbf{I}))^{1 / 2}\right)^{2}
$$

Proof. - Suppose $a^{\prime}>a>\mathrm{I}$. Let D denote the shell obtained on shrinking $\mathrm{D}_{\mathrm{G}}\left(a^{\prime}\right)$ by the factor $\frac{a}{a^{\prime}}$. Then, by Lemma (7.2), we get

$$
\bmod \mathrm{D}_{\mathrm{G}}\left(a^{\prime}\right)=\bmod \mathrm{D} \geqslant \bmod \mathrm{D}_{\frac{a}{a^{\prime}, 1}}+\bmod \mathrm{D}_{\mathrm{G}}(a)
$$

Thus

$$
\log \Phi_{n}\left(a^{\prime}\right) \geqslant \log \frac{a^{\prime}}{a}+\log \Phi_{n}(a)
$$

and this gives the first assertion.
To prove the second assertion, let D denote the shell obtained from $\mathrm{D}_{\mathrm{T}}(b)$ by translating $x \rightarrow\left(x_{1}+\mathrm{I}, x_{2}, \ldots, x_{n}\right)$ and then shrinking by the factor $(b+1)^{-1 / 2}$.

Let $\mathrm{D}^{\prime}=\mathrm{D}_{\mathrm{G}}\left((b+1)^{1 / 2}\right)$. Let $\sigma$ denote the Moebius inversion in the unit sphere and let $u$ denote an extremal function for the shell D . Then $\sigma$ interchanges the components of the complement of D and $\sigma \mathrm{D}=\mathrm{D}$. Since $\sigma$ preserves $\int_{\mathrm{D}}|\nabla u|^{n} d \mathrm{D}$, the function $u \circ \sigma$ is extremal and thus $u \circ \sigma=\mathrm{I}-u$. Since $\sigma(x)=x$ for $|x|=1$, we find $u(x)=\frac{1}{2}$ for $|x|=\mathrm{I}$. From this it follows that the restriction to $\mathrm{D}^{\prime}$ of $2\left(u-\frac{\mathrm{I}}{2}\right)$ is extremal for $\mathrm{D}^{\prime}$.

Hence

$$
\begin{aligned}
\mathscr{C}\left(\mathrm{D}_{\mathbf{T}}(b)\right) & =\int_{\mathbb{D}}|\nabla u|^{n} d x=\int_{\mathrm{D}^{\prime}}|\nabla u|^{n} d x+\int_{\sigma\left(\mathrm{D}^{\prime}\right)}|\nabla u|^{n} d x \\
& =2 \int_{\mathrm{D}^{\prime}}|\nabla u|^{n} d x=\frac{2}{2^{n}} \int_{\mathrm{D}^{\prime}}|\nabla \mathbf{2}((u-\mathrm{I}) / 2)|^{n} d x \\
& =2^{-(n-1)} \mathscr{C}\left(\mathrm{D}^{\prime}\right) .
\end{aligned}
$$

Thus

$$
\bmod \mathrm{D}_{\mathrm{T}}(b)=2 \bmod \mathrm{D}_{\mathrm{G}}\left((b+\mathrm{I})^{1 / 2}\right)
$$

or

$$
\Psi_{n}(b)=\left(\Phi_{n}\left((b+\mathrm{I})^{1 / 2}\right)\right)^{2}
$$

## § 8. SPHERICAL SYMMETRIZATION

The idea we are about to describe goes back to Jacob Steiner who defined for any plane region R and any line L in the plane, the symmetrization $\mathrm{R}^{*}$ as follows: At each point $x \in \mathrm{~L}$ such that the perpendicular line $\mathrm{S}_{x}$ to L at $x$ meets R , one places a line segment along $\mathrm{S}_{x}$ centered at $x$ whose length is the length of $\mathrm{S}_{x} \cap \mathrm{R}$; the union of these 82
segments is $R^{*}$. Schwarz introduced the symmetrization of a solid region in 3 -space with respect to a line L. If one thinks of Schwarz's planes perpendicular to L as a family of concentric spheres of infinite radius, then one is led to the definition of the spherical symmetrization of a set R in $\mathbf{R}^{n}$ with respect to the line $\mathrm{L}=\{(t, 0, \ldots, o) ;-\infty<t \leqslant \mathrm{o}\}$.

For each sphere $\mathrm{S}_{r}:|x|=r$ meeting the set R , place along $\mathrm{S}_{r}$, centered at $(-r, 0, \ldots, o)$ a spherical cap (of dimension $n-1$ ) of spherical measure equal to the measure of $S_{r} \cap R$; the union of these caps is denoted $R^{*}$.

To complete this definition, we have to specify whether the spherical caps are open or closed. One does the expected: If R is open, the spherical cap is taken open and includes the point $(r, o, \ldots, o)$ if and only if $S_{r} \subset R$. If R is closed, the spherical cap is taken closed.

It is clear that for any connected set E , open or closed, the set $\mathrm{E}^{*}$ is connected, and that $\mathrm{E}^{*}$ is open if E is open. Since the complement of $\mathrm{E}^{*}$ is clearly the spherical symmetrization of the complement of E with respect to the positive $x_{1}$-axis, it follows that $\mathrm{E}^{*}$ is closed if E is closed, and the complement of $\mathrm{E}^{*}$ is connected if the complement of E is connected. Also, $\mathrm{E}_{1}^{*} \subset \mathrm{E}_{2}^{*}$ if $\mathrm{E}_{1} \subset \mathrm{E}_{2}$.

Definition. - Let D be a shell in $\mathbf{R}^{n}$ and let $\mathrm{C}_{0}$ denote the bounded component of its complement. The sphere-symmetrization of D is the set $\mathrm{D}^{0}$ defined by

$$
\mathrm{D}^{0}=\left(\mathrm{D} \cup \mathrm{C}_{0}\right)^{*}-\mathrm{C}_{0}^{*} .
$$

It follows from the remarks above that $\mathrm{D}^{0}$ is a shell if D is a shell. For the Grötzsch and Teichmüller shells, we have $D_{G}^{0}=D_{G}, D_{T}^{0}=D_{T}$.

This section is devoted to the proof that $\mathscr{C}\left(\mathrm{D}^{0}\right) \leqslant \mathscr{C}(\mathrm{D})$ for any shell D , a fact which is suggested by isoperimetric considerations. In the course of comparing a shell with its symmetrization, we shall use the following notion.

Let $u$ be a real-valued continuous function on $\mathbf{R}^{n}$. For any real number $a$, set

$$
\begin{aligned}
\mathrm{F}_{a} & =u^{-1}(-\infty, a] \\
\mathrm{G}_{a} & =u^{-1}(-\infty, a) \\
\Delta_{a} & =u^{-1}(a)=\mathrm{F}_{a}-\mathrm{G}_{a} \\
\Delta_{a}^{*} & =\mathrm{F}_{a}^{*}-\mathrm{G}_{a}^{*} .
\end{aligned}
$$

Define the real valued function $u^{*}$ on $\mathbf{R}^{n}$ by:

$$
u^{*}(x)=a \text { if and only if } x \in \Delta_{a}^{*} .
$$

Clearly $\bigcap_{b>a} \mathrm{~F}_{b}=\mathrm{F}_{a}$, and $\bigcup_{b<a} \mathrm{~F}_{b}=\mathrm{G}_{a}$. Hence $\bigcap_{b>a} \mathrm{~F}_{b}^{*}=\mathrm{F}_{a}^{*}$ and $\bigcup_{b<a} \mathrm{~F}_{b}^{*}=\mathrm{G}_{a}^{*}$. It follows that $u^{*-1}(-\infty, a)=\mathrm{G}_{a}^{*}$ and is open. In addition, $u^{*-1}(-\infty, a]=\mathrm{F}_{a}^{*}$ is closed. Thus $u^{*-1}(b, a)$ is open for $b<a$ and therefore $u^{*}$ is continuous on $\mathbf{R}^{n}$. The function $u^{*}$ is called the spherical symmetrization of $u$. It has the property:

For any interval $\mathrm{I}=(a, b),[a, b]$, or $(a, b], u^{-1}(\mathrm{I}) \cap \mathrm{S}_{r}$ and $u^{*-1}(\mathrm{I}) \cap \mathrm{S}_{r}$ have the same measure.

We wish to compare $\mathscr{C}(\mathrm{D})$ and $\mathscr{C}\left(\mathrm{D}^{0}\right)$. This we do by comparing $\int_{\mathrm{D}}|\nabla u|^{n} d x$
and $\int_{\mathbb{D}^{0}}\left|\nabla u^{*}\right|^{n} d x$ for piecewise linear approximations $u$ to smoothly admissible functions. In the proof below we require the following two results:

The first result is an isoperimetric inequality:
Theorem A. - Let E be a subset of the simply connected space X of constant curvature, and let $\mathrm{T}_{s}(\mathrm{E})$ be a tubular neighborhood of radius $s$ (that is, the union of all balls of radius $s$ and center in E$)$. Then for fixed volume $m(\mathrm{E})$, the volume $m\left(\mathrm{~T}_{s}(\mathrm{E})\right.$ ) is minimized for E a ball.

This theorem, for convex bodies in euclidean space is known as the " BrunnMinkowski inequality ". It was first proved for arbitrary measurable sets in euclidean space by Lusternik, in C.R. Acad. Sci. U.R.S.S., v. 3 (1935), pp. 55-58. Theorem A is proved by E. Schmidt, in Math. Nachr., v. i (1948), pp. 8r-157. We shall require this theorem only for the cases $\mathrm{X}=\mathrm{S}^{n}$ and $\mathrm{X}=\mathbf{R}^{n}$.

The usual form of the isoperimetric theorem concerning the boundary $\partial \mathrm{E}(=\mathrm{E}-\mathrm{int} \mathrm{E})$ is a consequence of Theorem A. Namely, for any subset E of a metric space, one can define the $k$-dimensional Hausdorff measure

$$
m_{k}(\mathrm{E})=\lim _{s \rightarrow 0} \inf \sum_{\mathrm{V}} b_{k}(\text { radius } \mathrm{U})^{k}
$$

where $\{U\}$ is a countable covering of $E$ of balls $U$ having radius less than $s$, and

$$
b_{k}=2^{k} \Gamma\left(\frac{\mathrm{r}}{2}\right)^{k-1} \Gamma\left(\frac{k+\mathrm{I}}{2}\right) \Gamma(k+\mathrm{I})^{-1}
$$

that is, the volume of the unit ball in $\mathbf{R}^{k}$. In the highest dimension, Hausdorff measure coincides with Lebesgue measure on $\mathrm{S}^{n}$ and $\mathbf{R}^{n}$. Moreover if E is a closed region (i.e. $\mathrm{E}=\overline{\mathrm{int} \mathrm{E}}$ ) of either $\mathrm{S}^{n}$ of $\mathbf{R}^{n}$, with a piecewise smooth boundary, then

$$
m_{n-1}(\partial \mathrm{E})=\lim _{s \rightarrow 0} \frac{m\left(\mathrm{~T}_{s}(\mathrm{E})\right)-m(\mathrm{E})}{s} .
$$

Let $\hat{\mathrm{E}}$ denote the ball such that $m_{n}(\mathrm{E})=m_{n}(\hat{\mathrm{E}})$. (In the case of $\mathrm{S}^{n}$, the ball $\hat{\mathrm{E}}$ is a spherical cap of course). Then Theorem A implies

Theorem $A^{\prime}$. - $m_{n-1}(\partial \mathrm{E}) \geqslant m_{n-1}(\partial \hat{\mathrm{E}})$.
Notation. - Given $\mathrm{EcS}_{r}=\left\{x,|x|=r, x \in \mathbf{R}^{n}\right\}$ we shall be considering below the tubular neighborhoods $\mathrm{T}_{s}(\mathrm{E})$ in $\mathrm{S}_{r}$ and in $\mathbf{R}^{n}$. To avoid ambiguity, we shall denote by $\mathrm{T}_{s}(\mathrm{E})$ the tubular neighborhood in $\mathbf{R}^{n}$ and by $\widehat{\mathrm{T}}_{s}(\mathrm{E})$ the tubular neighborhood in $\mathrm{S}_{r}$.

The second theorem we shall refer to is:
Theorem B. - Let u be a real valued function on the domain $\mathrm{D} \subset \mathbf{R}^{n}$ satisfying the uniform Lipschitz condition

$$
|u(x)-u(y)|<\mathrm{M}|x-y| \quad \text { for } \quad x, y \in \mathrm{D} .
$$

Let $f$ be a continuous function on D . Then

$$
\int_{D} f|\nabla u| d x=\int_{-\infty}^{\infty}\left(\int_{u^{-1}(t)} f d \sigma\right) d t
$$

where d $\sigma$ denotes Hausdorff $(n-1)$-measure on $u^{-1}(t)$.

This theorem was first proved by E. di Giorgi (Ann. Mat. Pur. Appl., ser. 4, v. 36 (1954), pp. 191-213) under the hypothesis that $u$ is differentiable and was generalized independently by H. Federer and L. C. Young to Lipschitz maps $u$ of an $n$-manifold in to an $m$-manifold, $n \geqslant m$ (cf. H. Federer, Curvature measures, Trans. Amer. Math. Soc., v. 93 (1959), pp. $4^{\text {r8-431). }}$

Lemma (8.1). - Let F be a closed region in $\mathbf{R}^{n}$. Then

$$
\mathrm{T}_{s}(\mathrm{~F})^{*} \supset \mathrm{~T}_{s}\left(\mathrm{~F}^{*}\right) \quad \text { for any } \quad s>0
$$

where $\mathrm{T}_{s}$ denotes a tubular neighborhood of radius $s$.
Proof. - It suffices to show that for each $r>0$,

$$
\mathrm{S}_{r} \cap \mathrm{~T}_{s}(\mathrm{~F})^{*} \supset \mathrm{~S}_{r} \cap \mathrm{~T}_{s}\left(\mathrm{~F}^{*}\right) ;
$$

that is

$$
m_{n-1}\left(\mathrm{~S}_{r} \cap \mathrm{~T}_{s}(\mathrm{~F})\right) \geqslant m_{n-1}\left(\mathrm{~S}_{r} \cap \mathrm{~T}_{s}\left(\mathrm{~F}^{*}\right)\right)
$$

We have

$$
\mathrm{S}_{r} \cap \mathrm{~T}_{s}(\mathrm{~F})=\bigcup_{s \leqslant t \leqslant s} \mathrm{~S}_{r} \cap \mathrm{~T}_{s}\left(\mathrm{~S}_{r-t} \cap \mathrm{~F}\right)
$$

For any radius $r_{1}, r_{2}$, we have

$$
\mathrm{S}_{r_{2}} \cap \mathrm{~T}_{s}\left(\mathrm{~S}_{r_{1}} \cap \mathrm{~F}\right)=\hat{\mathrm{T}}_{\theta r_{2}}\left(\frac{r_{2}}{r_{1}}\left(\mathrm{~S}_{r_{1}} \cap \mathrm{~F}\right)\right)
$$

where $\theta$ is the angle defined by $s^{2}=r_{1}^{2}+r_{2}^{2}-2 r_{1} r_{2} \cos \theta$.
Hence, by Theorem A

$$
m_{n-1}\left(\mathrm{~S}_{r} \cap \mathrm{~T}_{s}\left(\mathrm{~S}_{r-t} \cap \mathrm{~F}\right)\right) \geqslant m_{n-1}\left(\mathrm{~S}_{r} \cap \mathrm{~T}_{s}\left(\mathrm{~S}_{r-t} \cap \mathrm{~F}^{*}\right)\right)
$$

Consequently,

$$
\begin{aligned}
m_{n-1}\left(\mathrm{~S}_{r} \cap \mathrm{~T}_{s}(\mathrm{~F})\right) & \geqslant \sup _{t} m_{n-1}\left(\mathrm{~S}_{r} \cap \mathrm{~T}_{s}\left(\mathrm{~S}_{r-t} \cap \mathrm{~F}\right)\right) \\
& \geqslant \sup _{t} m_{n-1}\left(\mathrm{~S}_{r} \cap \mathrm{~T}_{s}\left(\mathrm{~S}_{r-t} \cap \mathrm{~F}^{*}\right)\right) \\
& =m_{n-1}\left(\mathrm{U}_{-s \leqslant i \leqslant s} \mathrm{~S}_{r} \cap \mathrm{~T}_{s}\left(\mathrm{~S}_{r-t} \cap \mathrm{~F}^{*}\right)\right)
\end{aligned}
$$

since $\mathrm{T}_{s}\left(\mathrm{~F}^{*}\right)$ is a solid of revolution. Hence

$$
m_{n-1}\left(\mathrm{~S}_{r} \cap \mathrm{~T}_{s}(\mathrm{~F})\right) \geqslant m_{n-1}\left(\mathrm{~S}_{r} \cap \mathrm{~T}_{s}\left(\mathrm{~F}^{*}\right)\right)
$$

Lemma (8.2). - Let F be a closed region in $\mathbf{R}^{n}$ with a piecewise smooth boundary. Then

$$
m_{n}\left(\mathrm{~F}^{*}\right)=m_{n}(\mathbf{F})
$$

and

$$
m_{n-1}\left(\partial \mathbf{F}^{*}\right) \leqslant m_{n-1}(\partial \mathbf{F}) .
$$

Proof. - Set $\mathrm{V}(r)=m_{n}\left(\mathrm{~F} \cap \mathrm{~B}_{r}\right), \mathrm{V}^{*}(r)=m_{n}\left(\mathrm{~F}^{*} \cap \mathrm{~B}_{r}\right)$.
Then $\frac{d \mathrm{~V}}{d r}=m_{n-1}\left(\mathrm{~F} \cap \mathrm{~S}_{r}\right)=m_{n-1}\left(\mathrm{~F}^{*} \cap \mathrm{~S}_{r}\right)=\frac{d \mathrm{~V}^{*}}{d r}$. It follows that $\mathrm{V}(r)=\mathrm{V}^{*}(r)$ and $m_{n}\left(\mathbf{F}^{*}\right)=m_{n}(\mathbf{F})$.

To prove the second assertion, we observe

$$
\begin{aligned}
m_{n-1}(\partial \mathbf{F}) & =\lim _{s \rightarrow 0} \frac{m\left(\mathrm{~T}_{s}(\mathbf{F})\right)-m(\mathbf{F})}{s} \\
& =\lim _{s \rightarrow 0} \frac{m\left(\mathrm{~T}_{s}(\mathbf{F})^{*}\right)-m\left(\mathbf{F}^{*}\right)}{s} \\
& \geqslant \lim _{s \rightarrow 0} \frac{m\left(\mathrm{~T}_{s}\left(\mathbf{F}^{*}\right)\right)-m\left(\mathbf{F}^{*}\right)}{s}, \quad \text { by Lemma (8.1). }
\end{aligned}
$$

Hence $m_{n-1}(\partial \mathbf{F}) \geqslant m_{n-1}\left(\partial F^{*}\right)$.
Lemma (8.3). - Let u be a real-valued function on $\mathrm{R}^{n}$ satisfying the uniform Lipschtiz condition

$$
|u(x)-u(y)| \leqslant \mathrm{M}|x-y| \quad \text { for } \quad x, y \in \mathbf{R}^{n} .
$$

Let $u^{*}$ denote the spherical symmetrization of $u$. Then

$$
\left|u^{*}(x)-u^{*}(y)\right| \leqslant \mathrm{M}|x-y| \quad \text { for } \quad x, y \in \mathbf{R}^{n} .
$$

Proof. - Let $\mathrm{F}_{i}$ denote the set in $\mathbf{R}^{n}$ on which $u(x) \leqslant t$. It suffices to prove that if $x \in \mathrm{~T}_{s}\left(\mathrm{~F}_{t}^{*}\right)$, then $u^{*}(x) \leqslant t+\mathrm{Ms}$; that is, $\mathrm{T}_{s}\left(\mathrm{~F}_{t}^{*}\right) \subset \mathrm{F}_{t+\mathrm{M}_{s}}^{*}$. By Lemma (8. I) and our hypothesis

$$
\mathrm{T}_{s}\left(\mathrm{~F}_{t}^{*}\right) \subset \mathrm{T}_{s}\left(\mathrm{~F}_{t}\right)^{*} \subset \mathrm{~F}_{t+\mathrm{M}_{s}}^{*}
$$

Lemma (8.4). - Let F be a fnite union of simplexes in $\mathbf{R}^{n}$, and let $u$ be a continuous piecewise linear function on $\mathbf{R}^{n}$, linear on each simplex of F and whose gradient vanishes outside F . Let $u^{*}$ denote the spherical symmetrization of $u$. Let $\varphi$ be a piecerwise continuous function on $\mathbf{R}^{2}$. Set $\Delta_{t}=u^{-1}(t), \Delta_{t}^{*}=u^{*-1}(t)$. Then

$$
\int_{\mathbf{R}^{n}} \varphi(|x|, u(x)) d x=\int_{\mathbf{R}^{n}} \varphi\left(|x|, u^{*}(x)\right) d x
$$

and for all but a finite number of $t$,

$$
\int_{\Delta_{t}} \varphi(|x|, t) d \Delta_{t} \geqslant \int_{\Delta_{t}^{*}} \varphi(|x|, t) d \Delta_{t}^{*}
$$

Proof. - Set $\mathrm{W}(r)=\int_{\mathrm{B}_{r}} \varphi(|x|, u(x)) d x, \mathrm{~W}^{*}(r)=\int_{\mathrm{B}_{r}} \varphi\left(|x|, u^{*}(x)\right) d x$.
Then $\frac{d \mathrm{~W}}{d r}=\int_{\mathrm{S}_{r}} \varphi(r, u(x)) d \mathrm{~S}_{r}$ and a similar formula holds for $\frac{d \mathrm{~W}^{*}}{d r}$.
From the definition of $u^{*}$ it follows at once that the set of points on $\mathrm{S}_{r}$ for which $a<u \leqslant b$ and $a<u^{*} \leqslant b$ respectively have the same measure, $a$ and $b$ being arbitrary. It follows at once that $\frac{d \mathrm{~W}}{d r}=\frac{d \mathrm{~W}^{*}}{d r}$ and $\mathrm{W}(r)=\mathrm{W}^{*}(r)$. This establishes the first assertion.

Turning to the second assertion, let $t_{1}, \ldots, t_{p}$ be the values of $u$ on the vertices of F . Then for $t$ distinct from these values, $\Delta_{t}=u^{-1}(t)$ is the boundary of $\mathrm{F}_{t}=u^{-1}[-\infty, t]$; that is, $\mathrm{F}_{t}=\overline{\operatorname{Int} \mathrm{F}_{t}}$. Moreover $\Delta_{t}$ is piecewise smooth and $\Delta_{t} \cap \mathrm{~S}_{r}$ is the boundary of $\mathrm{F}_{t} \cap \mathrm{~S}_{r}$ for each $r>0$, provided $t \neq t_{1}, \ldots, t_{p}$.

Let $\mathrm{F}_{t} ; r_{1}, r_{2}=\left(\mathrm{F}_{t} \cap \mathrm{~B}_{r_{\mathbf{2}}}\right)-\left(\mathrm{F}_{t} \cap \operatorname{Int} \mathrm{~B}_{r_{2}}\right)$, where $\mathrm{o}<r_{1}<r_{2}$, and let

$$
\Delta_{t} ; r_{1}, r_{2}=\Delta_{t} \cap\left(\mathrm{~F}_{t} ; r_{1}, r_{2}\right)
$$

Similarly define $\mathrm{F}_{t}^{*} ; r_{1}, r_{2}$ and $\Delta_{t}^{*} ; r_{1}, r_{2}$. Then

$$
\begin{aligned}
& \partial \mathrm{F}_{t} ; r_{1}, r_{2}=\left(\Delta_{t} ; r_{1}, r_{2}\right) \cup\left(\mathrm{F}_{t} \cap \mathrm{~S}_{r_{1}}\right) \cup\left(\mathrm{F}_{t} \cap \mathrm{~S}_{r_{2}}\right) \\
& \partial \mathrm{F}_{t}^{*} ; r_{1}, r_{2}=\left(\Delta_{t} ; r_{1}, r_{2}\right) \cup\left(\mathrm{F}_{t}^{*} \cap \mathrm{~S}_{r_{1}}\right) \cup\left(\mathrm{F}_{t}^{*} \cap \mathrm{~S}_{r_{2}}\right)
\end{aligned}
$$

By definition $m_{n-1}\left(\mathrm{~F}_{t} \cap \mathrm{~S}_{r_{i}}\right)=m_{n-1}\left(\mathrm{~F}_{t}^{*} \cap \mathrm{~S}_{r_{i}}\right)(i=\mathrm{I}, 2)$. Hence by Lemma (8.2)

$$
m_{n-1}\left(\Delta_{t} ; r_{1}, r_{2}\right) \geqslant m_{n-1}\left(\Delta_{t}^{*} ; r_{1}, r_{2}\right)
$$

for $\mathrm{o}<r_{1}<r_{2}$. Let M denote the set function on the positive $r$-axis given by $\mathrm{M}\left(r_{1}, r_{2}\right)=m_{n-1}\left(\Delta_{t} ; r_{1}, r_{2}\right)$, and similarly set $\mathrm{M}^{*}\left(r_{1}, r_{2}\right)=m_{n-1}\left(\Delta_{t}^{*} ; r_{1}, r_{2}\right)$. Then

$$
\int_{\Delta_{t}} \varphi(|x|, t) d \Delta_{t}=\int_{r>0} \varphi(r, t) d \mathrm{M} \geqslant \int_{r>0} \varphi(r, t) d \mathrm{M}^{*}=\int_{\Delta_{t}^{*}} \varphi(|x|, t) d \Delta_{t}^{*}
$$

Theorem (8.1). - Let D be a shell in $\mathrm{R}^{n}$ and let $\mathrm{D}^{0}$ be its sphere-symmetrization. Then

$$
\mathscr{C}(\mathrm{D}) \geqslant \mathscr{C}\left(\mathrm{D}^{0}\right)
$$

Proof. - Let $u$ be a smoothly admissible function for D. By the method of § 3, one can approximate D by a piecewise linear function $u^{\prime}$ and the function $u^{\prime}$ can be so chosen that
I. $u^{\prime}$ is admissible for D and $\nabla u^{\prime}$ has as support a finite union of simplexes.
2. $\left|\nabla u-\nabla u^{\prime}\right|$ is arbitrarily small.

Since the spherical symmetrization of $u^{\prime}$ will then be admissible for $\mathrm{D}^{0}$, it suffices, to prove Theorem (8.I), to show

$$
\int_{\mathbf{R}^{n}}|\nabla u|^{n} d x \geqslant \int_{\mathbf{R}^{n}}\left|\nabla u^{*}\right|^{n} d x
$$

for any piecewise linear function which " lives" on a finite union of simplexes. Let $t_{1}<t_{2}<\ldots<t_{h}$ denote the values of $u$ on the vertices of these simplexes.

Set $f^{*}=\left|\nabla u^{*}\right|^{n-1}$. We claim $f^{*}(x)=\varphi\left(|x|, u^{*}(x)\right)$ with $\varphi$ a piecewise continuous function on $\mathbf{R}^{2}$. For by the axial symmetry of $u^{*}, \nabla u$ is determined by its restriction to the upper half ( $x_{1}, x_{2}$ )-coordinate plane P ; let $\mathrm{G}_{i}$ denote the subset of points in $\mathbf{R}^{n}$ for which $t_{i}<u^{*}(x)<t_{i+1}$. Then the map of $\mathbf{R}^{n}, \pi: x \rightarrow\left(|x|, u^{*}(x)\right)$ gives a homeomorphism $\pi_{i}$ of $\mathrm{G} \cap \mathrm{P}_{i}$. For by definition of spherical symmetrization, $u^{*}$ is a monotone function on the semi-circle $S_{r} \cap P$ and it is not constant on any open subset of $S_{r} \cap G_{i} \cap P$. Since in addition $u^{*}$ is piecewise differentiable, it is strictly monotone on the arc $\mathrm{S}_{r} \cap \mathrm{G}_{i} \cap \mathrm{P}$ and provides a homeomorphism of the arc for each $r$. It follows that $\pi_{i}$, the restriction of $\pi$ to $G_{i} \cap \mathrm{P}$, is a homeomorphism. Therefore the restriction of $f^{*}$ to $\mathrm{G}_{i}$ can be expressed by $f^{*}(x)=\left(f^{*} \circ \pi_{i}^{-1}\right)(\pi(x))$. Set $\varphi=f^{*} \circ \pi_{i}^{-1}$ on $\pi\left(\mathrm{G}_{i}\right)(i=1, \ldots, h-\mathrm{I})$. Define $\varphi$ to be zero on the complement of $\pi\left(\mathbf{R}^{n}\right)$. It is clear that $\varphi$ is piecewise continuous on $\mathbf{R}^{2}$.

$$
\begin{aligned}
\text { We have } \\
\begin{array}{rlrl}
\int_{\mathbf{R}^{n}}\left|\nabla u^{*}\right|^{n} d x & =\int_{\mathbf{R}^{n}}\left|\nabla u^{*}\right|^{n-1}\left|\nabla u^{*}\right| d x=\int f^{*}\left|\nabla u^{*}\right| d x & \\
& =\int_{0}^{1}\left(\int_{u^{*-1}(t)} f^{*} d \sigma\right) d t, & & \\
& \leqslant \int_{0}^{1}\left(\int_{u^{-1}(t)} f d \sigma\right) d t, \quad \text { where } f(x)=\varphi(|x|, u(x)), & & \text { by Lemma (8.4) } \\
& =\int_{\mathbf{R}^{n}} f|\nabla u| d x, & & \text { by Theorem B } \\
& \leqslant\left(\int_{\mathbf{R}^{n}} f^{\frac{n}{n-1}}\right)^{\frac{n-1}{n}}\left(\int_{\mathbf{R}^{n}}|\nabla u|^{n}\right)^{\frac{1}{n},} & & \text { by Theorem B } \\
& =\left(\int_{\mathbf{R}^{n}}\left|\nabla u^{*}\right|^{n}\right)^{\frac{n-1}{n}}\left(\int_{\mathbf{R}^{n}}|\nabla u|^{n}\right)^{\frac{1}{n},} & & \text { by Hölder's inequality }
\end{array}
\end{aligned}
$$

Hence

$$
\int_{\mathbf{R}^{n}}|\nabla u|^{*} d x \leqslant \int_{\mathbf{R}^{n}}|\nabla u|^{n} .
$$

Theorem (8.2). - Let D be a shell in $\mathbf{R}^{n}$, and let $\mathrm{C}_{0}, \mathrm{C}_{1}$ denote the bounded and unbounded components of its complement. Let $x_{0} \in \mathrm{C}_{0}$, and set $a=\sup _{y \in \mathrm{C}_{0}} d\left(x_{0}, y\right), b=\inf _{y \in \mathrm{C}_{1}} d\left(x_{0}, y\right)$.
Then

Then

$$
\bmod \mathrm{D} \leqslant \log \Psi_{n}(b / a)
$$

Proof. - Let $\mathrm{D}^{0}$ denote the sphere symmetrization of D based on $x_{0}$ as the origin and the line L as axis. Then clearly $\mathrm{C}_{0}^{*} \supset[-a, o]$, and $\mathrm{C}_{1}^{*} \supset[b, \infty]$. Hence $\mathrm{D}^{0}$ separates the boundaries of the shell $\mathrm{D}^{\prime}=\mathbf{R}^{n}-[-a, 0]-[b, \infty]$ and we have

$$
\bmod \mathrm{D} \leqslant \bmod \mathrm{D}^{0} \leqslant \bmod \mathrm{D}^{\prime}=\bmod \mathrm{D}_{\mathbf{T}}(b / a)=\log \Psi_{n}(b / a)
$$

## § 9. QUASI CONFORMAL MAPPINGS

In this section we shall define quasi-conformal mappings and establish the characteristic properties that will be needed in § to to prove:

A quasiconformal mapping of the open n-ball onto itself extends to a homeomorphism of the boundary, and this boundary homeomorphism is quasi-conformal.

Let $\varphi: \mathrm{D} \rightarrow \mathrm{D}^{\prime}$ be a homeomorphism of a Riemannian manifold D onto the Riemannian manifold $\mathrm{D}^{\prime}$. We associate with $\varphi$ the functions $\mathrm{H}_{\varphi}, \mathrm{I}_{\varphi}, \mathrm{J}_{\varphi}, \mathrm{L}_{\varphi}$, and $l_{\varphi}$ defined as follows:

$$
\begin{aligned}
\mathrm{L}_{\varphi}(p, r) & =\sup _{d(q, p)=r} d(\varphi(q), \varphi(p)) \\
l_{\varphi}(p, r) & =\inf _{d(q, p)=r} d(\varphi(q), \varphi(p)) \\
\mathrm{H}_{\varphi}(p) & =\varlimsup_{r \rightarrow 0} \frac{\mathrm{~L}(p, r)}{l(p, r)} \\
\mathrm{I}_{\varphi}(p) & =\varlimsup_{r \rightarrow 0} \frac{\mathrm{~L}(p, r)}{r} \\
\mathrm{~J}_{\varphi}(p) & =\varlimsup_{r \rightarrow 0} \frac{m\left(\varphi\left(\mathrm{~T}_{r}(p)\right)\right)}{m\left(\mathrm{~T}_{r}(p)\right)}
\end{aligned}
$$

Here we have written $\mathrm{L}=\mathrm{L}_{\varphi}, l=l_{\varphi}$, and similarly we write $\mathrm{H}, \mathrm{I}$, J when the homeomorphism $\varphi$ is unambiguous in the context. We note that the tubular neighborhood $\mathrm{T}_{r}(p)$ of the point $p$ is merely an open ball.

Some elementary observations should be noted.

Properties (1), (2), and (3) are obvious. The first property in (4) comes from

$$
\left(\frac{\mathrm{L}_{\varphi}(p, t)}{t}\right)^{n}=\left(\frac{\mathrm{L}_{\varphi}(p, t)}{l_{\varphi}(p, t)}\right)^{n} \cdot\left(\frac{l_{\varphi}(p, t)}{t}\right)^{n} .
$$

The second property in (4) comes from setting $t=\mathrm{L}_{\varphi-1}(\varphi(p), s)$ and observing, by (3), that

$$
\left(\frac{t}{s}\right)^{n}=\left(\frac{t}{l_{\varphi}(p, t)}\right)^{n}=\left(\frac{\mathrm{L}_{\varphi}(p, t)}{l_{\varphi}(p, t)}\right)^{n} \cdot\left(\frac{t}{\mathrm{~L}_{\varphi}(p, t)}\right)
$$

and

$$
\frac{t}{\mathrm{~L}_{\varphi}(p, t)}=\frac{l_{\varphi-1}(p, r)}{r}
$$

where $r=\mathrm{L}_{\varphi}(p, t)$.
We say that a mapping $\varphi$ is differentiable at a point $p$ if it has a differential at $p$. If the domain and image of $\varphi$ are in $\mathbf{R}^{n}$, the differential $\dot{\varphi}$ of $\varphi$ at $p$ is a linear mapping such that

$$
\varphi(x)=\varphi(p)+\dot{\varphi}(x-p)+\varepsilon(x, p)
$$

where $\lim _{x \rightarrow p} \frac{|\varepsilon(x, p)|}{|x-p|}=0$. If $\varphi$ is differentiable at $p$, then the above functions satisfy $\mathrm{H}_{\varphi}(p)=\mathrm{H}_{\varphi}(\mathrm{o})$ if $\varphi \neq \mathrm{o}, \mathrm{I}_{\varphi}(p)=\mathrm{I}_{\dot{\varphi}}(\mathrm{o}), \mathrm{J}_{\varphi}(p)=\mathrm{J}_{\varphi}(\mathrm{o})$.

For a linear mapping $\varphi: \mathbf{R}^{n} \rightarrow \mathbf{R}^{n}$, we have

$$
\mathrm{L}_{\varphi}(0, r)^{2}=\sup _{|x|=r}\langle\varphi x, \varphi x\rangle=\sup _{|x|=r}\left\langle{ }^{t} \varphi \varphi x, x\right\rangle=\lambda_{1}^{2} r^{2}
$$

where $\lambda_{1}^{2}$ is the maximum eigenvalue of the positive definite self-adjoint operator ${ }^{t} \varphi \varphi$. Similarly $l_{\varphi}(0, r)^{2}=\lambda_{n}^{2} r^{2}$ where $\lambda_{n}^{2}$ is the smallest eigenvalue of ${ }^{t} \varphi \varphi$. Since

$$
J_{\varphi}(0)=|\operatorname{det} \varphi|=\left(\operatorname{det}^{\prime} \varphi \varphi\right)^{1 / 2}=\lambda_{1} \lambda_{2} \ldots \lambda_{n},
$$

where $\lambda_{1}^{2} \geqslant 0$ and $\lambda_{1}^{2} \geqslant \lambda_{2}^{2} \geqslant \ldots \geqslant \lambda_{n}^{2}$ are the eigenvalues of ${ }^{t} \varphi \varphi$, we have

$$
\lambda_{1}^{n} \leqslant\left(\lambda_{1} / \lambda_{n}\right)^{n-1} \lambda_{1} \lambda_{2} \ldots \lambda_{n} .
$$

It follows that
(5) If $\varphi$ is differentiable at $p$, then $\mathrm{I}_{\varphi}(p)^{n} \leqslant \mathrm{H}_{\varphi}(p)^{n-1} \mathrm{~J}_{\varphi}(p)$ since $\mathrm{I}_{\varphi}(p)=\mathrm{J}_{\varphi}(p)=0$ if $\dot{\varphi}=0$.

We shall require the following result, known as the Rademacher-Stepanoff Theorem.
Let $f$ be a measurable function on a measurable set E such that

$$
\lim _{\Delta x \rightarrow 0}\left|\frac{f(x+\Delta x)-f(x)}{\Delta x}\right|<\infty \quad \text { a.e. }
$$

Then $f$ is differentiable a.e. in E .
This result was proved by H. Rademacher (Math. Annalen, v. 79, pp. 340-359 (1919)) under the stronger hypothesis that $\mathrm{I}_{j}<\infty$ everywhere in the domain E . The result quoted here is given by W. Stepanoff, Rec. Math. Soc. Math., Moscou, v. 32, pp. 51I-526 (1925), cf. also Saks, p. 3 Io.

Definition. - A homeomorphism $\varphi$ of a Riemannian manifold D onto a Riemannian manifold $\mathrm{D}^{\prime}$ is called quasi-conformal if $\mathrm{H}_{\varphi}$ is bounded on D . The mapping $\varphi$ is called K -quasi-conformal if $\mathrm{H}_{\varphi}$ is bounded on D and $\mathrm{H}_{\varphi}(x) \leqslant \mathrm{K}$ almost everywhere on D .

For example, a conformal mapping is I-quasi-conformal.
Theorem (9.1). - A quasi-conformal mapping is differentiable almost everywhere.
Proof. - By definition, $\mathrm{J}_{\varphi}(p)$ is the upper derivative of the set function $\mathrm{E} \rightarrow m(\varphi(\mathrm{E}))$. Hence by Lebesgue's Theorem, $\mathrm{J}_{\varphi}<\infty$ a.e. (Saks, p. II5). Hence by (4), $\mathrm{I}_{\varphi}<\infty$ a.e. Theorem (9.I) now follows from the Rademacher-Stepanoff Theorem.

Let $\mathrm{E} \subset \mathbf{R}^{n}$ and let $a$ be a positive real number. Let $\{U\}$ be a denumerable cover of E by balls of radius not exceeding $a$. We write

$$
\Lambda(\mathrm{E}, a)=\inf _{\{\mathrm{U}\}}\left(\sum_{\mathrm{U}} 2 \operatorname{rad} \mathrm{U}\right)
$$

Then $\lim _{a \rightarrow 0} \Lambda(\mathrm{E}, a)=m_{1}(\mathrm{E})$, the I-dimensional Hausdorff measure of the set E .
Lemma (9.1). - Let $\varphi$ be a topological mapping of a domain $\mathrm{D} \subset \mathbf{R}^{n}$ onto a domain in $\mathbf{R}^{n}, n>\mathrm{I}$, such that $\mathrm{H}_{\varphi}(x)<\mathrm{K}$ for all $x \in \mathrm{D}-\mathrm{P}$, where P is a hyperplane in $\mathbf{R}^{n}$. Then $\varphi$ is ACL in D and $\varphi^{-1}$ is ACL in $\varphi(\mathrm{D})$.

Proof. - The lemma is equivalent to the assertion that $\varphi$ is ACL in D if either (i) $\mathrm{H}_{\varphi}$ is bounded in $\mathrm{D}-\mathrm{P}$ or (ii) $\mathrm{H}_{\varphi-1}$ is bounded in $\varphi(\mathrm{D}-\mathrm{P})$.

To prove that $\varphi$ is ACL in D, it suffices to prove that for almost all closed line segments X interior to D and parallel to the $x_{1}$ axis, $\varphi$ is absolutely continuous on X ; that is, $m_{1}(\varphi(E))=0$ if $E$ is a subset of $X$ with $m_{1}(E)=0$. Set

$$
\tau(\mathrm{X})=\varlimsup_{t \rightarrow 0} \frac{m\left(\varphi\left(\mathrm{~T}_{t}(\mathrm{X})\right)\right)}{t^{n-1}}
$$

Then by Lebesgue's Theorem on differentiation of set functions (Saks, Theory of the Integral, p. II5), $\tau(\mathrm{X})<\infty$ for almost all line segments X . The lemma will be established once we have proved for any compact set $\mathrm{E} \subset \mathrm{X}-\mathrm{P}$ (which we can assume to be X minus one point)

$$
\begin{equation*}
m_{1}(\varphi(\mathrm{E}))^{n} \leqslant \mathrm{~A} \tau(\mathrm{X}) m_{1}(\mathrm{E})^{n-1} \tag{6}
\end{equation*}
$$

where $A$ is a constant depending on the bound for $H_{\varphi}$ in case (i) and the bound for $H_{\varphi-1}$ in case (ii).

First we take up case (i); that is we assume $\mathrm{H}(x)<\mathrm{K}$ for all $x \in \mathrm{D}-\mathrm{P}$. We can assume without loss of generality that there is a positive number $b$ satisfying for all $p \in \mathrm{E}$

$$
\begin{equation*}
\mathrm{L}(p, t) \leqslant \mathrm{K} l(p, t) \quad \text { whenever } \quad o<t<b \tag{7}
\end{equation*}
$$

For let $\mathrm{E}_{b}$ be the set of $p \in \mathrm{E}$ such that (7) is satisfied. Since $\mathrm{E}_{b}$ is compact and expands to E as $b \rightarrow 0$, the validity of (6) for E follows from its validity for $\mathrm{E}_{b}$. Next we fix a positive number $a$ (which ultimately will tend to zero) and we select a positive $c$ such that $\mathrm{L}(p, c)<a$ for all $p \in \mathrm{E}$; this we can do by the uniform continuity of $\varphi$ on E . We can find a real number $t, 0<t<\inf (b, c)$ and points $p_{1}, p_{2}, \ldots, p_{\mathrm{N}}$ in the linear set E such that the balls $\mathrm{T}_{i}\left(p_{k}\right)$ cover $\mathrm{E}, \quad\left|p_{i}-p_{j}\right| \geqslant(|i-j|-\mathrm{I}) t$ for all $i, j=\mathrm{I}, \ldots, \mathrm{N}$, and $\mathrm{N} t \leqslant m_{1}(\mathbf{E})+a$. Namely, choose $t<\min (b, c)$ so that $m_{1}(\mathrm{~T}(\mathrm{E}) \cap \mathrm{X})<m_{1}(\mathrm{E})+a$, divide the line containing E into equal disjoint half-open intervals of length $t$, select in each interval meeting E a point of E , and arrange these points $p_{1}, p_{2}, \ldots, p_{\mathrm{N}}$ in order. We note in passing that for any real number K , no point of the ball $\left|x-p_{k}\right|<\mathrm{K} t$ lies in more than 3 K of the balls with center at $p_{1}, \ldots, p_{\mathrm{N}}$ and radius $\mathrm{K} t$. Set

$$
s_{k}=\mathrm{L}\left(p_{k}, t\right) \quad(k=\mathrm{I}, \ldots, \mathrm{~N}) .
$$

Clearly $s_{k}<a$.
From (7) we get

$$
s_{k} / \mathrm{K} \leqslant l\left(\boldsymbol{p}_{k}, t\right)
$$

and thus by (1) and (2)

$$
\mathrm{T}_{s_{k} / K}\left(\varphi\left(p_{k}\right)\right) \subset \varphi\left(\mathrm{T}_{t}\left(p_{k}\right)\right) \subset \mathrm{T}_{s_{k}}\left(\varphi\left(p_{k}\right)\right)
$$

Set $\quad \mathrm{B}_{k}=\mathrm{T}_{s_{k} / \mathrm{K}}\left(\varphi\left(p_{k}\right)\right), \mathrm{B}^{k}=\mathrm{T}_{s_{k}}\left(\varphi\left(p_{k}\right)\right)$.
Since the $\mathrm{B}^{k}$ form a cover of $\varphi(\mathrm{E})$ by balls of radius less than $a$, we have

$$
\Lambda(\varphi(\mathrm{E}), a) \leqslant \sum_{k} 2 s_{k} .
$$

By Hölder's inequality,

$$
\begin{equation*}
\Lambda(\varphi(\mathrm{E}), a)^{n} \leqslant 2^{n} \mathbf{N}^{n-1} \sum_{k} s_{k}=b_{n}^{-1} 2^{n} \mathbf{N}^{n-1} \sum_{k} m\left(\mathbf{B}^{k}\right) \tag{8}
\end{equation*}
$$

where $b_{n}$ denotes the volume of the unit $n$-ball. Since $\mathrm{B}_{k} \subset \varphi\left(\mathrm{~T}_{t}\left(p_{k}\right)\right)$, each point of $\mathrm{B}_{k}$ lies in at most three of $\mathrm{B}_{1}, \ldots, \mathrm{~B}_{\mathrm{N}}$. Hence

$$
\sum_{k} m\left(\mathbf{B}_{k}\right) \leqslant 3 m\left(\bigcup_{k} \mathrm{~B}_{k}\right)
$$

and

$$
\sum_{k} m\left(\mathrm{~B}^{k}\right)=\mathrm{K}^{n} \sum_{k} m\left(\mathrm{~B}_{k}\right) \leqslant 3 \mathrm{~K}^{n} m\left(\bigcup_{k} \mathrm{~B}_{k}\right) \leqslant 3 \mathrm{~K}^{n} m\left(\varphi\left(\mathrm{~T}_{i}(\mathrm{E})\right)\right) .
$$

Consequently

$$
\Lambda(\varphi(\mathrm{E}), a)^{n} \leqslant b_{n}^{-1} 2^{n} \cdot 3 \mathrm{~K}^{n}(\mathrm{~N} t)^{n-1} \frac{m\left(\varphi\left(\mathrm{~T}_{t}(\mathrm{E})\right)\right)}{t^{n-1}} \leqslant \mathrm{~A}\left(m_{1}(\mathrm{E})+a\right)^{n-1} \frac{m\left(\varphi\left(\mathrm{~T}_{t}(\mathrm{X})\right)\right)}{t^{n-1}}
$$

Letting $a \rightarrow 0$, we find formula (6).

In the case (ii), we can assume instead of (7), the condition

$$
\mathrm{L}_{\varphi-1}(\varphi(p), t) \leqslant \mathrm{K} l_{\varphi-1}(\varphi(p), t) \quad 0<t<b
$$

for all $p \in \mathrm{E}$. Again we fix $a$, and select $c$ such that $\mathrm{L}(p, c)<a$. Again we choose $t<\inf (b, c)$ and the points $p_{1}, \ldots, p_{\mathrm{N}}$ so that the balls $\mathrm{T}_{t}\left(p_{k}\right)$ cover E ,

$$
\left|p_{i}-p_{j}\right| \geqslant(|i-j|-1) t
$$

for all $i, j$, and $\mathrm{N} t \leqslant m_{1}(\mathrm{E})+a$.
Again we set $s_{k}=\mathrm{L}_{\varphi}\left(p_{k}, t\right)(k=\mathrm{I}, \ldots, \mathrm{N})$. Then

$$
l_{\varphi-1}\left(\varphi\left(p_{k}\right), s_{k}\right)=t, \quad \text { and } \quad \mathrm{L}_{\varphi-1}\left(p_{k}, s_{k}\right)<\mathrm{K} t
$$

by (3) and ( $7^{\prime}$ ). Hence

$$
\varphi\left(\mathrm{T}_{i}\left(p_{k}\right)\right) \subset \mathrm{T}_{s_{k}}\left(\varphi\left(p_{k}\right)\right) \subset \varphi\left(\mathrm{T}_{\mathrm{K} i}\left(p_{k}\right)\right)
$$

As before, formula (8) applies. But in this case we estimate $\sum_{k} m\left(\mathbf{B}^{k}\right)$ by observing that each point of $\mathrm{B}^{k}$ meets at most 3 K of $\varphi\left(\mathrm{T}_{\mathrm{K} t}\left(p_{k}\right)\right), k=\mathrm{I}, \ldots, \mathrm{N}$.

Consequently, $\sum_{k} m\left(\mathrm{~T}_{s_{k}}\left(\varphi\left(p_{k}\right)\right)\right) \leqslant 3 \mathrm{~K} m\left(\mathrm{U}_{k} \mathrm{~T}_{s_{k}}\left(\varphi\left(p_{k}\right)\right)\right)$, and

$$
\Lambda(\varphi(\mathrm{E}), a)^{n} \leqslant b_{n}^{-2} 2^{n} 3 \mathrm{KN}^{n-1} m\left(\varphi\left(\mathrm{~T}_{\mathrm{K} t}(\mathrm{X})\right)\right) .
$$

Therefore

$$
m_{1}(\varphi(\mathrm{E})) \leqslant b_{n}^{-1} 2^{n} 3 \mathrm{~K}^{n}\left(m_{1}(\mathrm{E})\right)^{n-1} \tau(\mathrm{X})=\mathrm{A} m_{1}(\mathrm{E})^{n-1} \tau(\mathrm{X}) .
$$

The proof of Lemma (9.I) is now complete.
Lemma (9.2). - Let $\varphi: \mathbf{R} \rightarrow \mathbf{R}^{\prime}$ be a homeomorphism of domains in $\mathbf{R}^{n}$. Assume $\varphi$ is ACL in R and $\mathrm{I}_{\varphi}^{n} \leqslant \mathrm{~K}^{n-1} \mathrm{~J}_{\varphi}$ a.e. in $\mathrm{R}, \mathrm{K}$ being a constant. Then

$$
\bmod \varphi(\mathrm{D}) \leqslant \mathrm{K} \bmod \mathrm{D}
$$

for all shells D in R .
Proof. - Let $u^{\prime}$ be a differentiable admissible function on the shell $\varphi(\mathrm{D})$ and set $u=u^{\prime} \circ \varphi$. Then $u$ is ACL on D and hence admissible on D .

Now

$$
\begin{aligned}
|\nabla u|(p) & =\varlimsup_{q \rightarrow p} \frac{|u(q)--u(p)|}{|q-p|} \\
& \leqslant \varlimsup_{q \rightarrow p} \frac{|u(q)-u(p)|}{|\varphi(q)-\varphi(p)|} \cdot \varlimsup_{q \rightarrow p} \frac{|\varphi(q)-\varphi(p)|}{|q-p|}
\end{aligned}
$$

(9)

$$
|\nabla u|(p) \leqslant\left|\nabla u^{\prime}\right|(\varphi(p)) \cdot \mathrm{I}(p) .
$$

Hence

$$
|\nabla u|^{n} \leqslant\left(\left|\nabla u^{\prime}\right| \circ \varphi\right) \cdot \mathrm{I}^{n} \leqslant \mathrm{~K}^{n-1}\left(\left|\nabla u^{\prime}\right|^{n} \circ \varphi\right) \cdot \mathrm{J}
$$

and

$$
\mathscr{C}(\mathrm{D}) \leqslant \mathrm{K}^{n-1} \int_{\mathrm{D}}\left(\left|\nabla u^{\prime}\right|^{n} \circ \varphi\right) \cdot \mathrm{J} d x \leqslant \mathrm{~K}^{n-1} \int_{\varphi(\mathrm{D})}\left|\nabla u^{\prime}\right|^{n} d x
$$

for all differentiable admissible functions $u^{\prime}$ on $\varphi(\mathrm{D})$. It follows that $\mathscr{C}(\mathrm{D}) \leqslant \mathrm{K}^{n-1} \mathscr{C}(\varphi(\mathrm{D}))$ and $\bmod \varphi(\mathrm{D}) \leqslant K \bmod D$.

Note. - It follows from (5), Theorem (9.1), and Lemma (9.1) that $\bmod \varphi(\mathrm{D}) \leqslant \mathrm{K} \bmod \mathrm{D}$ if $\varphi$ is K-quasi-conformal. It follows from (4) and Lemma (9. I) that $\bmod \varphi^{-1}(\mathrm{D}) \leqslant \mathrm{K}^{\frac{n}{n-1}} \bmod \mathrm{D}$ if $\varphi$ is K -quasi-conformal. Following Lemma (9.3) below, we will know that $\varphi^{-1}$ is quasi-conformal, and being differentiable with non-zero differential almost everywhere, we will be able to assert that $\varphi^{-1}$ is also K -quasi-conformal. It will therefore follow that $\bmod \varphi^{-1}(\mathrm{D}) \leqslant \mathrm{K} \bmod \mathrm{D}$.

Theorem (9.2). - Let $\varphi$ be a homeomorphism of a domain R in $\mathbf{R}^{n}$ onto a domain in $\mathbf{R}^{n}$. Assume $n>I_{1}$ and

$$
\bmod \varphi(\mathrm{D}) \leqslant \mathrm{K} \bmod \mathrm{D}
$$

for all shells D in R . Then $\varphi$ is quasi-conformal and $\mathrm{H}_{\varphi}<\Psi_{n}(\mathrm{I})^{\mathrm{K}}$ on R .
Proof. - We must prove that $\mathrm{H}_{\varphi}$ is bounded on R .
Let $p \in \mathrm{R}$, and let $r$ be any positive number such that the ball $\mathrm{T}_{r}(p) \subset \mathrm{R}$. Then

$$
\varphi^{-1}\left(\mathrm{~T}_{\langle(p, r)}(\varphi(p))\right) \subset \mathrm{T}_{r}(p) \subset \varphi^{-1}\left(\mathrm{~T}_{\mathrm{L}(p, r)}(\varphi(p))\right) .
$$

Let D denote the shell whose complement consists of

$$
\begin{aligned}
& \mathrm{C}_{0}=\text { closure of } \varphi^{-1}\left(\mathrm{~T}_{l(p, r)}(\varphi(p))\right) \\
& \mathrm{C}_{1}=\text { exterior of } \varphi^{-1}\left(\mathrm{~T}_{\mathrm{L}(p, r)}(\varphi(p))\right) .
\end{aligned}
$$

Then the spherical symmetrization $\mathrm{D}^{0}$ separates the boundaries of the Teichmüller-like shell $\mathrm{D}^{\prime}$ whose complement consists of the intervals $\mathrm{C}_{0}^{\prime}=[-r, o], \mathrm{C}_{1}^{\prime}=[r, \infty]$ along the axis of the symmetrization. Hence

$$
\bmod \varphi(\mathrm{D}) \leqslant \mathrm{K} \bmod \mathrm{D} \leqslant \mathrm{~K} \bmod \mathrm{D}^{0} \leqslant \mathrm{~K} \bmod \mathrm{D}_{\mathrm{T}}(\mathrm{I})=\mathrm{K} \log \Psi_{n}(\mathrm{I}) .
$$

But $\varphi(\mathrm{D})$ is the spherical shell $\mathrm{D}_{l, \mathrm{~L}}$ so that we get

$$
\log \frac{\mathrm{L}(p, r)}{l(p, r)} \leqslant \mathrm{K} \log \Psi_{n}(\mathrm{I}) .
$$

It follows at once that $\mathrm{H}_{\varphi}(p) \leqslant \Psi_{n}(\mathrm{I})^{\mathrm{K}}$ for all $p \in \mathrm{R}$. Hence $\varphi$ is quasi-conformal.
We may summarize our conclusions in the next theorem :
Theorem (9.3). - The following assertions are equivalent for a homeomorphism $\varphi$ of a domain R in $\mathbf{R}^{n}$ onto a domain of $\mathbf{R}^{n}, n>_{\mathrm{I}}$ :
a) $\mathrm{H}_{\varphi}$ is bounded on R ; that is $\varphi$ is quasi-conformal.
b) $\varphi$ is ACL in R and $\mathrm{I}^{n}<\mathrm{KJ}$ a.e. for some constant K .
c) $\bmod \varphi(\mathrm{D}) \leqslant \mathrm{K} \bmod \mathrm{D}$ for all shells D in R for some constant K .
d) $\varphi^{-1}$ is quasi-conformal.
e) $\mathrm{K}^{-1} \bmod \mathrm{D} \leqslant \bmod \varphi(\mathrm{D})$ for all shells D in R for some constant K .

Proof. - $a) \Rightarrow b) \Rightarrow c) \Rightarrow a$ ) has already been proved. a) $\Rightarrow d$ ) follows from (4), Lemma (9.1), and $b$ ), $d) \Rightarrow a$ ) is obviously equivalent to $a) \Rightarrow d$ ), and $d) \Leftrightarrow e$ ) by $a$ ) and $c$ ).

Theorem (9.4). - Let $\varphi$ be a quasi-conformal mapping of a domain R of $\mathbf{R}^{n}$ onto a domain of $\mathbf{R}^{n}, n>\mathrm{I}$. Then $\mathrm{J}_{\varphi}>0$ a.e. in R , the set function $\mathrm{E} \rightarrow m(\varphi(\mathrm{E}))$ is absolutely continuous in R and

$$
m(\varphi(\mathrm{E}))=\int_{\mathrm{E}} \mathrm{~J} d x
$$

for any measurable set E in R .
Proof. - Let Z be the set of points in R at which $\mathrm{J}_{\varphi}$ is zero. To prove that the set $Z$ has measure zero, it suffices to prove that $Z$ has no points of density; that is, a point $p$ such that $\varlimsup_{r \rightarrow 0} \frac{m\left(\mathrm{~T}_{r}(p) \cap \mathrm{Z}\right)}{m\left(\mathrm{~T}_{r}(p)\right)}=\mathrm{I}$. (Saks, p. 129).

Let D denote the spherical shell $a<|x-p|<2 a$ centered at $p$ and let $u^{\prime}$ be a differentiable admissible function on $\varphi(\mathrm{D})$. Set $u=u^{\prime} \circ \varphi$. Then $u$ is an admissible function on D . Let $y$ denote a variable point in the coordinate plane $x=x_{1}(p)$, and let $\mathrm{X}_{y}$ denote the intersection with D of the line through D parallel to the $x_{1}$-axis. Then we have

$$
\int_{\mathrm{D}}|\nabla u| \geqslant \int_{|q-p|<a} \int_{\mathrm{X}_{y}}|\nabla u| d x_{1} d y \geqslant 2 b_{n-1} a^{n-1}
$$

where $b_{n-1}$ is the volume of the unit ( $n-1$ )-ball.
On the other hand, since $|\nabla u|(p) \leqslant\left|\nabla u^{\prime}\right|(\varphi(p)) \cdot \mathbf{I}(p)$,

$$
2 b_{n-1} a^{n-1} \leqslant \int_{\mathrm{D}}\left(\left|\nabla u^{\prime}\right| \circ \varphi\right) . \mathrm{I} d x .
$$

Assume $\varphi$ to be K-quasi-conformal, we have $I^{n} \leqslant \mathrm{~K}^{n-1} \mathrm{~J}$, and applying Hölder's inequality,

$$
\begin{aligned}
\left(2 b_{n-1} a^{n-1}\right)^{n} \leqslant & \left(\int_{\mathrm{D}-\mathrm{Z}}\left(\left|\nabla u^{\prime}\right| \circ \varphi\right) \mathrm{I} d x\right)^{n} \leqslant \mathrm{~K}^{n-1} m(\mathrm{D}-\mathrm{Z})^{n-1} \int_{\mathrm{D}-\mathrm{Z}}\left(\left|\nabla u^{\prime}\right|^{n} \circ \varphi\right) \cdot \mathrm{J} d x \\
& \leqslant \mathrm{~K}^{n-1} m(\mathrm{D}-\mathrm{Z})^{n-1} \int_{\varphi(\mathrm{D})}\left|\nabla u^{\prime}\right|^{n} d x
\end{aligned}
$$

It follows that $\left(2 b_{n-1} a^{n-1}\right)^{n} \leqslant \mathrm{~K}^{n-1} m(\mathrm{D}-\mathrm{Z})^{n-1} \mathscr{C}(\varphi(\mathrm{D}))$ or

$$
m(\mathrm{D}-\mathrm{Z}) \geqslant \mathrm{K}^{-1} \mathrm{~A} a^{n} \bmod \varphi(\mathrm{D})
$$

where A is a constant depending only on the dimension $n$. Now

$$
\bmod \varphi(\mathrm{D}) \geqslant \mathrm{K}^{-n /(n-1)} \bmod \mathrm{D}
$$

by Lemma (9.2) applied to $\varphi^{-1}$. Since $\bmod \mathbf{D}=\bmod D_{a, 2 a}=\log 2$, we find

$$
\frac{m(\mathrm{D}-\mathrm{Z})}{m\left(\mathrm{~T}_{2 a}(p)\right)}>\mathrm{A}^{\prime}>\mathrm{o}
$$

where $\mathrm{A}^{\prime}$ is a constant depending only on K and $n$. It follows that $p$ is not a point of density of $Z$, and that $Z$ has measure zero.

Let P denote the set of points $p$ in R at which $\mathrm{J}_{\varphi}(p)=\infty$. Then any point of $\varphi(\mathrm{P})$ at which $\varphi^{-1}$ is differentiable must lie in $Z_{\varphi-1}$, the set of zeros of $J_{\varphi-1}$. Since $\varphi^{-1}$ is quasi-conformal, we find that $\varphi(\mathrm{P})$ is a set of measure zero. It follows from the De la

Vallée Poussin Decomposition Theorem (Saks, p. 127) that $\mathrm{E} \rightarrow m(\varphi(\mathrm{E})$ ) is absolutely continuous and $m(\varphi(\mathrm{E}))=\int_{\mathrm{E}} \mathrm{J} d x$.

Corollary. - If $\varphi$ is K -quasi-conformal, then $\varphi^{-1}$ is K -quasi-conformal.
Proof. - If $n=\mathrm{I}$, the assertion is obvious. We can assume therefore that $n>_{\mathrm{I}}$. We can assert now that $\dot{\varphi} \neq 0$ a.e. and thus $\mathrm{H}_{\varphi}(p)=\mathrm{H}_{\dot{\varphi}_{p}}$ for almost all $p$ if $\varphi$ is quasiconformal, and similarly for $\varphi^{-1}$. The Corollary follows from the fact that $H_{\varphi}=H_{\varphi-1}$ for a non-singular linear mapping $\varphi$.

Remark. - As a consequence of (5) and Lemma (9.2), we have for any shell D

$$
\mathrm{K}^{-1} \bmod \mathrm{D} \leqslant \bmod \varphi(\mathrm{D}) \leqslant \mathrm{K} \bmod \mathrm{D}
$$

if $\varphi$ is K -quasi-conformal on D .

## § 10. THE BOUNDARY MAP OF THE $n$-BALL

Theorem (10.1). - Let $\varphi$ be a quasi-conformal mapping of an open ball in $\mathbf{R}^{n}$ onto itself. Then $\varphi$ extends to a homeomorphism of the closed $n$-ball.

Proof. - Mapping the domain of $\varphi$ onto the upper half space

$$
\mathrm{X}=\left\{\left(x_{1}, \ldots, x_{n}\right) ; x_{n}>0\right\}
$$

via a Moebius transformation, the theorem is seen to be equivalent to the assertion: a quasi-conformal mapping $\varphi: \mathrm{X} \rightarrow \mathrm{Y}=\{y ;|y|<\mathrm{I}\}$ extends to a continuous mapping at any point $x$ of the boundary of $X$. For convenience, we take $x=0$.

The proof is by contradiction. If $\lim _{p \rightarrow 0} \varphi(p)(p \in \mathrm{X})$ does not exist, we can find two sequences $\left\{p_{k}\right\}$ and $\left\{q_{k}\right\}$ in X approaching o with $\lim _{k \rightarrow \infty} \varphi\left(p_{k}\right)=p^{\prime}, \lim _{k \rightarrow \infty} \varphi\left(q_{k}\right)=q^{\prime}$ and $\left|q^{\prime}-p^{\prime}\right|=a>o$. Denoting by $\overline{p q}$ the line segment joining two points $p$ and $q$, we select points $p_{0}^{\prime}$ and $q_{0}^{\prime}$ in Y such that $d\left(\overline{p_{0}^{\prime} p_{k}^{\prime}}, \overline{q_{0}^{\prime} q_{k}}\right)>a$ for all large $k$, where $p_{k}^{\prime}=\varphi\left(p_{k}\right)$, $q_{k}^{\prime}=\varphi\left(q_{k}\right)$. Set $p_{0}=\varphi^{-1}\left(p_{0}^{\prime}\right), q_{0}=\varphi^{-1}\left(q_{0}^{\prime}\right)$. Then for $\sup \left(\left|p_{k}\right|,\left|q_{k}\right|\right)<r<\inf \left(\left|p_{0}\right|,\left|q_{0}\right|\right)$, the hemisphere $\mathrm{S}_{r}^{+}=\left\{x ;|x|=r, x_{n}>0\right\}$ meets the curves $\varphi^{-1}\left(p_{0}^{\prime} p_{k}^{\prime}\right)$ and $\varphi^{-1}\left(\overline{q_{0}^{\prime} q_{k}^{\prime}}\right)$. For each such $r$ at least one of the coordinate functions of $\varphi(x)=\left(\varphi_{1}(x), \ldots, \varphi_{n}(x)\right)$ satisfies

$$
\underset{S_{n}^{+}}{\operatorname{osc}} \varphi_{i}>a / \sqrt{n}
$$

Hence

$$
\sum_{i} \int_{0}^{\infty}\left(\underset{\mathrm{s}_{r}^{+}}{\operatorname{osc}} \varphi_{i}\right)^{n} \frac{d r}{r}=\infty
$$

By Lemma (9. I), $\varphi_{i}$ is ACL in X. Applying Lemma (4.4), we get for each $i=1, \ldots, n$

$$
\begin{aligned}
\int_{0}^{\infty}\left(\underset{S_{r}}{\operatorname{osc}} \varphi_{i}\right)^{n} \frac{d r}{r} & \leqslant 2 \mathrm{~A} \int_{\mathrm{X}}\left|\nabla \varphi_{i}\right|^{n} d x \leqslant 2 \mathrm{~A} \int_{\mathrm{X}} \mathrm{I}_{\varphi}^{n} d x \leqslant 2 \mathrm{~A} \int_{\mathrm{X}} \mathrm{~K}^{n-1} \mathrm{~J}_{\varphi} d x \leqslant 2 \mathrm{AK}^{n-1} \int_{\mathrm{Y}} d y \\
& <\infty
\end{aligned}
$$

This yields a contradiction.

Theorem (10.2). - Let $\varphi$ be a K-quasi-conformal mapping of an open ball in $\mathbf{R}^{n}$ onto itself, $n \geqslant 2$, and let $\bar{\varphi}$ denote the boundary homeomorphism induced by $\varphi$. Then $\bar{\varphi}$ is $\bar{\Psi}_{n}(\mathrm{I})^{\mathrm{K}}$-quasiconformal.

Proof. - Denote the extended homeomorphism to the closed ball by $\varphi$ and set $\mathrm{K}^{\prime}=\bar{\Psi}_{n}(\mathrm{I})^{\mathrm{K}}$. We wish to show that $\mathrm{H}_{\bar{\varphi}}<\mathrm{K}^{\prime}$ on the boundary ( $n-\mathrm{r}$ )-sphere. Given any point $p$ on the boundary, no generality is lost in assuming that $\bar{\varphi}\left(p_{\infty}\right)=p_{\infty}$ for some point $p_{\infty}$ on the boundary distinct from $p$ : we need merely compose $\varphi$ with some rotations of the closed ball to achieve this. Let $\theta$ be a Moebius transformation of $\mathbf{R}^{n} \cup \infty$ carrying the closed ball onto the upper half space $x_{n} \geqslant 0$ with $p_{\infty}$ going into $\infty$. Set $\psi=\theta \varphi \theta^{-1}$. Then $\mathrm{H}_{\psi}(x)=\mathrm{H}_{\varphi}\left(\theta^{-1}(x)\right)$ for all $x$ in the open upper half space and $\mathrm{H}_{\bar{\psi}}(x)=\mathrm{H}_{\bar{\varphi}}\left(\theta^{-1} x\right)$ for all $x$ in the coordinate plane $\mathrm{P}: x_{n}=0$.

We extend $\psi$ to a homeomorphism of $\mathbf{R}^{n}$ to $\mathbf{R}^{n}$, defining

$$
\psi\left(x_{1}, \ldots, x_{n-1},-x_{n}\right)=-\psi\left(x_{1}, \ldots, x_{n}\right) .
$$

Then $\mathrm{H}_{\psi}(x)<\mathrm{K}$ for all $x \in \mathbf{R}^{n}-\mathrm{P}$. It follows from Lemma (9.1), that $\varphi$ is ACL in $\mathbf{R}^{n}$ and it follows from $\S 9$, formula (5), that $\mathrm{I}^{n} \leqslant \mathrm{~K}^{n-1} \mathrm{~J}$ a.e. in $\mathbf{R}^{n}$. Hence by Lemma (9.2),

$$
\bmod \psi(\mathrm{D}) \leqslant K \bmod \mathrm{D}
$$

for all shells D in $\mathbf{R}^{n}$. Applying Theorem (9.2), we get $\mathrm{H}_{\psi}(x) \leqslant \mathrm{K}^{\prime}$ for all $x \in \mathbf{R}^{n}$. In particular

$$
\mathrm{H}_{\bar{\varphi}}(p)=\mathrm{H}_{\bar{\psi}}\left(\theta^{-1}(p)\right) \leqslant \mathrm{H}_{\psi}(\theta(p)) \leqslant \mathrm{K}^{\prime}
$$

and thus $\bar{\varphi}$ is $\mathrm{K}^{\prime}$-quasi-conformal.

## § II. AN ERGODIC THEOREM

Lemma (1x. $\mathbf{1}$ ). - Let S be a non-abelian group of the form B.A where B is a normal subgroup isomorphic to the additive group of real or complex numbers and A is a subgroup isomorphic to the multiplicative group of positive real numbers, and with multiplication

$$
(b \cdot a) \circ\left(b^{\prime} \cdot a^{\prime}\right)=\left(b+a b^{\prime}\right) \cdot a a^{\prime} .
$$

Let $\rho$ be a continuous unitary representation of S on a hilbert space V and let $v$ be a vector in V such that $\rho(\mathrm{A}) v=v$. Then $\rho(\mathrm{S}) v=v$.

Proof. - For any $a \in \mathrm{~A}$ and $b \in \mathrm{~B}$ we have $a \circ b \circ a^{-1}=a b$ and thus

$$
\begin{aligned}
\langle\rho(b) v, v\rangle & =\left\langle\rho\left(a^{n}\right) \rho(b) v, \rho\left(a^{n}\right) v\right\rangle \\
& =\left\langle\rho\left(a^{n} \circ b \circ a^{-n}\right) \rho\left(a^{n}\right) v, \rho\left(a^{n}\right) v\right\rangle \\
& =\left\langle\rho\left(a^{n} b\right) v, v\right\rangle
\end{aligned}
$$

for all $n>0$. Choosing $|a|<1$, we get as $n \rightarrow \infty$

$$
\langle\rho(b) v, v\rangle=\langle\rho(\text { identity }) v, v\rangle=\langle v, v\rangle .
$$

By Schwarz's inequality, $\rho(b) v=v$ for all $b \in \mathrm{~B}$ and hence $\rho(\mathrm{S}) v=v$.

Note. - The foregoing lemma was first proved by Mautner by expressing any unitary representation as a direct integral of irreducible ones and verifying the assertion for the irreducible representations of S . The proof given here is taken from an article of L. Auslander and L. Green in Amer. 7. Math., v. 88 (1966), p. 58, who in turn attribute the underlying idea to J. von Neumann and I. E. Segal. It may be further noted that $p(B) v=v$ if we make the weaker hypothesis that $v$ is an eigenvector for A.

Lemma (1I.2). - Let W be a measure space with finite total measure $\mu$. Let A be a group of measurable measure-preserving transformations of W such that the only A-invariant functions in $\mathrm{L}^{2}(\mathrm{~W}, \mu)$ are the constant functions. Let $\mathrm{A}^{+}$be a semi-group in A such that $\mathrm{A}^{+}$and its inverse generate A. Assume that there is a separable topology on A with respect to which the representation of A on $\mathrm{L}^{2}(\mathrm{~W}, \mu)$ is continuous. Let $\left\{\mathrm{W}_{n}\right\}$ be a denumerable family of measurable sets in W . Then for almost all $x \in \mathrm{~W}, \mathrm{~A}^{+} x$ meets $\mathrm{W}_{n}$.

Proof. - Set $\mathrm{A}^{-}=\left(\mathrm{A}^{+}\right)^{-1}$. Consider the characteristic function $f$ of the measurable set $\mathrm{DW}_{n}$, where D is the semi-group generated by a denumerable dense subset of $\mathrm{A}^{-}$. For any $d \in \mathrm{D}$ we have $d \mathrm{DW}_{n} \subset \mathrm{DW}_{n}$ and since the action of $d$ is measurepreserving, $\mu\left(d \mathrm{DW}_{n}\right)=\mu\left(\mathrm{DW}_{n}\right)$. Hence $d f=f$ for all $d \in \mathrm{D}$. Since D is dense in $\mathrm{A}^{-}$, we find $\mathrm{A}^{-} f=f$ and therefore $\mathrm{A} f=f$. By hypothesis $f$ is constant almost everywhere and therefore $\mu\left(\mathrm{DW}_{n}\right)=\mu(\mathrm{W})$. Consequently $\mathrm{A}^{-} \mathrm{W}_{n}$ differs from W in merely a set of measure zero and $\mu\left(\bigcap_{n} \mathrm{~A}^{-} \mathrm{W}\right)=\mu(\mathrm{W})$. For any $x \in \bigcap_{n} \mathrm{~A}^{-} \mathrm{W}_{n}$, we find $\mathrm{A}^{+} x \cap \mathrm{~W}_{n}$ non-empty for all $n$.

Theorem (11.x). - Let G be a separable locally compact group and $\Gamma$ a discrete subgroup such that $\mathrm{G} / \Gamma$ has finite Haar measure. Let A be a subgroup of G isomorphic to the multiplicative group of positive real numbers and let $\mathrm{A}^{+}$be the semi-group in A given by $\{a ; a \leqslant \mathrm{I}\}$. Assume that $\mathbf{G}$ is generated by $\mathbf{A}$ and the family of subgroups isomorphic to the real numbers $\{\mathbf{B}\}$ such that B.A satisfies the multiplication rule: $(b . a) \circ\left(b^{\prime} \cdot a^{\prime}\right)=\left(b+a b^{\prime}\right) . a a^{\prime}$. Then for almost all cosets $x \Gamma, \mathrm{~A}^{+} x \Gamma$ is topologically dense in G .

Proof. - Consider the canonical representation of G on $\mathrm{L}^{2}(\mathrm{G} / \Gamma)$ and let $f$ be a function fixed under A. By Lemma (II, I), $f$ is fixed under B.A for all subgroups in the family $\{B\}$. By hypothesis, $\{B A\}$ generates $G$. Hence $f$ is fixed by $G$; that is, $f$ is a constant function on $G / \Gamma$. The hypotheses of Lemma (II.2) are satisfied by A and the conclusion of Theorem (II.I) follows by taking for $\left\{\mathrm{W}_{n}\right\}$ a denumerable base of open sets in $G / \Gamma$.

Note. - The subset AГ may well be closed.

## § 12. PROOF OF THE MAIN THEOREM

Theorem (12.1). - Let $G$ be the group of isometries of hyperbolic $n$-space $\mathrm{X} . \quad$ Let $\Gamma$ and $\Gamma^{\prime}$ be subgroups of G such that $\Gamma \backslash \mathrm{G}$ and $\Gamma^{\prime} \backslash \mathrm{G}$ have finite Haar measure. Let $\theta: \Gamma \rightarrow \Gamma^{\prime}$ be an isomorphism and $\varphi: \mathrm{X} \rightarrow \mathrm{X}$ a quasi-conformal homeomorphism such that

$$
\varphi(\gamma x)=\theta(\gamma) \varphi(x)
$$

for all $\gamma \in \Gamma, x \in \mathrm{X}$. Then $\theta$ extends to an inner automorphism of G provided $n>2$.

The proof is divided into several parts. We begin with some remarks on the algebraic structure of $G$.

1. We have defined hyperbolic $n$-space at the end of Section I as the unit ball $\mathrm{B}^{n}:|x|<\mathrm{I}$ in $\mathbf{R}^{n}$ with metric $\left(\mathrm{I}-|x|^{2}\right)^{-2} \sum_{i=1}^{n} d x_{i}^{2}$ and we have seen in Theorem (1.2) that the subgroup $\mathrm{G}^{\prime}$ of the Moebius group $\mathrm{GM}(n)$ which stabilizes a hemisphere in $\mathrm{S}^{n+1}$ operates on hyperbolic $n$-space transitively with isotropy subgroup $\mathrm{O}(n)$. It follows at once that any isometry of hyperbolic space differs from an isometry in $\mathrm{G}^{\prime}$ by an isometry T which keeps a point fixed and induces on the tangent space at some point $p$ the identity transformation. In such circumstances T leaves fixed all the geodesics through $p$ and is therefore the identity transformation. Hence the group $G$ of all isometries of $X$ is precisely the group $\mathrm{G}^{\prime}=\mathrm{GM}(n-\mathrm{I})=\mathrm{O}(\mathrm{I}, n) /( \pm \mathrm{I})$. In other words, when we take $\mathrm{B}^{n}=\left\{x ;|x|<_{\mathrm{I}}\right\}$ as our model of hyperbolic $n$-space X , each isometry of X induces a Moebius transformation on the boundary sphere $\mathrm{S}^{n-1}=\{x ;|x|=\mathrm{I}\}$ and conversely, each Moebius transformation of $\mathrm{S}^{n-1}$ extends to a unique isometry of X.

Let $p_{0}$ and $p_{\infty}$ be antipodal points on the $n-\mathrm{I}$ sphere $\mathrm{S}^{n-1}$ in $\mathbf{R}^{n}$. Let $\pi_{\infty}$ denote stereographic projection of $\mathrm{S}^{n-1}$ from $p_{\infty}$ onto the tangent plane $\mathrm{E}_{0}$ to $\mathrm{S}^{n-1}$ at $p_{0}$ and let $\mathrm{T}_{0}$ denote the ( $n-\mathrm{I}$ )-dimensional vector group of translations in the plane $\mathrm{E}_{0}$. Set $\mathrm{N}_{+}=\pi_{\infty}^{-1} \mathrm{~T}_{0} \pi_{\infty}$. Then $\mathrm{N}_{+} \subset \mathrm{G}$. Similarly, performing stereographic projection from $p_{0}$, we define the vector subgroup $\mathrm{N}_{-}$of $\mathrm{G}_{0}$. Let A denote the one-parameter subgroup of G corresponding to the homotheties $x \rightarrow t x$ where $t \in \mathbf{R}$ and $x \in \mathrm{E}_{0}$ which is regarded as a vector space with origin $p_{0}$. It is clear that $A$ corresponds equally to the homotheties from $p_{\infty}$ in the plane $E_{\infty}$. Then $A N_{+}$is a solvable subgroup of $G$ and similarly $A N_{-}$ is a solvable subgroup.

We have seen in Section I that $\mathrm{GM}(n-\mathrm{I})$ is generated by the reflection $\sigma$ in the equatorial plane bisecting $p_{0} p_{\infty}$ together with $\mathrm{AN}_{+}$. Clearly $\sigma^{2}=1, \sigma \mathrm{~A} \sigma=\mathrm{A}$ and $\sigma \mathrm{N}_{+} \sigma=\mathrm{N}_{-}$. Hence $\mathrm{G}=\mathrm{GM}(n-\mathrm{I})$ has exactly two connected components and $\mathrm{G}_{0}$, the connected component of the identity element in $G$, is generated by the solvable subgroups $A N_{+}$and $A N_{-}$. Let $B$ denote a one parameter vector subgroup of either $N_{+}$ or $N_{-}$. Then BA is a solvable subgroup of $G_{0}$; if we identify $B$ with the additive group of real numbers and $A$ with the multiplicative group of positive real numbers, we have for $b, b^{\prime} \in \mathrm{B}$ and $a, a^{\prime} \in \mathrm{A}$ :

$$
(b \cdot a) \circ\left(b^{\prime} \cdot a^{\prime}\right)=\left(b+a b^{\prime}\right) \cdot a a^{\prime}
$$

2. The fixed point set of the subgroup $A$ is the set of points $\left\{p_{0}, p_{\infty}\right\}$ and therefore the stabilizer of $\left\{p_{0}, p_{\infty}\right\}$ in G is $\mathrm{N}(\mathrm{A})$, the normalizer of $A$. Hence we may identify the space of cosets $G / N(A)$ with the orbit under $G$ of the $\mid$ set $\left\{p_{0}, p_{\infty}\right\}$, that is, the set of unordered pairs of distinct points $\{p, q\}$ of $S^{n-1}$. On the other hand each coset of $G / N(A)$ determines a unique conjugate of A ; namely $x \mathrm{~N}(\mathrm{~A}) \leftrightarrow x \mathrm{~A} x^{\prime}$. Given the distinct points $p, q$ in $\mathrm{S}^{n-1}$, we denote by $\mathrm{A}_{p, q}$ the conjugate of A corresponding to $\{p, q\} ; \mathrm{A}_{p, q}$ is the unique conjugate of A leaving fixed $p$ and $q$.

The hypotheses of Theorem (II.I) are satisfied by the triple $\mathrm{G}_{0}, \mathrm{\Gamma}_{\mathrm{n}} \mathrm{G}_{0}, \mathrm{~A}$. It follows that for almost all cosets $\Gamma x$ in the group $\Gamma G_{0}$ (which may be either $G_{0}$ or $G$ ), the subset $\Gamma x \mathrm{~A}^{+}$is topologically dense in $\Gamma \mathrm{G}_{0}$, where $\mathrm{A}^{+}$denotes the semi-group in A corresponding to $\{a ; a<\mathrm{I}\}$. Hence for almost all $x \in \mathrm{G}_{0}, \Gamma x \mathrm{~A}^{+} x^{-1}$ is topologically dense in $\Gamma \mathrm{G}_{0}$. In other words $\Gamma \mathrm{A}_{p, q}^{+}$is topologically dense in $\Gamma \mathrm{G}_{0}$ for almost all $\{p, q\}$.
3. By hypothesis we have a quasi-conformal mapping $\varphi: \mathrm{X} \rightarrow \mathrm{X}$ of hyperbolic $n$-spaces such that $\varphi(\gamma x)=\theta(\gamma) \varphi(x)$, for all $\gamma \in \Gamma, x \in \mathrm{X}$. The mapping $\varphi$ is thus a homeomorphism of the unit $n$-ball $\mathbf{B}^{n}$ onto itself which is quasi-conformal with respect to the hyperbolic metric $\left(\mathrm{I}-|x|^{2}\right)^{-2} \sum_{i=1}^{n} d x_{i}^{2}$. It follows at once that $\varphi$ is quasi-conformal with respect to the euclidean metric $\sum_{i=1}^{n} d x_{i}^{2}$. By Theorems (10.1) and (10.2) we may assert that $\varphi$ extends to a homeomorphism $\psi$ of the boundary $\mathrm{S}^{n-1}$ and $\psi$ is quasiconformal. Clearly

$$
\psi(\gamma p)=\theta(\gamma) \psi(p)
$$

for all $\gamma \in \Gamma, p \in \mathrm{~S}^{n-1}$. We shall prove that $\psi$ is a Moebius transformation if $n>2$.
4. Since $\psi: \mathrm{S}^{n-1} \rightarrow \mathrm{~S}^{n-1}$ is quasi-conformal, its differential $\dot{\psi}_{p}$ exists at almost all points $p \in \mathrm{~S}^{n-1}$, by Theorem (9.1). In addition $\Gamma \mathrm{A}_{p, q}^{+}$is topologically dense in $\Gamma \mathrm{G}_{0}$ for almost all $\{p, q\}$, when we identify the unordered pair of distinct point $\{p, q\}$ with an element of $\mathrm{G} / \mathrm{N}(\mathrm{A})$ and introduce a quotient measure on $\mathrm{G} / \mathrm{N}(\mathrm{A})$. The centralizer $\mathrm{Z}(\mathrm{A})$ is the fixer of $p_{0}$ and $p_{\infty}$. It is easy to verify that, up to a constant factor, the invariant quotient measure of $\mathrm{G} / \mathrm{Z}(\mathrm{A})=\left(\mathrm{S}^{n-1} \times \mathrm{S}^{n-1}\right.$-diagonal) correspond to the product of the standard measures on $\mathrm{S}^{n-1} \times \mathrm{S}^{n-1}$. Applying Theorem (III), we find for almost all points $p \in \mathrm{~S}^{n-1}$, there is a point $q$ distinct from $p$ such that $\Gamma \mathrm{A}_{p, q}^{+}$is topologically dense in $\Gamma G_{0}$. Finally, if $n>_{\mathrm{I}}, \dot{\psi}_{p}$ is invertible for almost all $p$, by Theorem (9.4).
5. Let $p$ be a point of $\mathrm{S}^{n-1}$ such that
a) the differential $\dot{\psi}_{p}$ exists and is invertible.
b) $\Gamma_{0} \mathrm{~A}_{p, q}^{+}$is topologically dense in $\mathrm{G}_{0}$ for some $q$.

We prove under these hypotheses that $\dot{\psi}_{p}$ is a conformal linear mapping of the tangent space to $\mathrm{S}^{n-1}$ at $p$. To reach this conclusion about $\dot{\psi}_{p}$, we may compose $\psi$ with a Moebius transformation of the closed unit ball and thus no generality is lost in assuming $\psi(p)=p$, and $\psi(q)=q$. Next we carry the pair $(p, q)$ into the antipodal points ( $p_{0}, p_{\infty}$ ) by an element $\tau \in G$ and we replace $\Gamma$ by $\tau \Gamma \tau^{-1}, \varphi$ by $\varphi \circ \tau^{-1}, \theta(\gamma)$ by $\theta\left(\tau^{-1} \gamma \tau\right)$. Thereby we arrive at the situation : $(p, q)=\left(p_{0}, p_{\infty}\right), \psi\left(p_{0}\right)=p_{0}$, and $\Gamma_{0} \mathrm{~A}_{p_{0}, p_{\infty}}^{+}$is topologically dense in $\mathrm{G}_{0}$.

We identify the tangent space to $\mathrm{S}^{n-1}$ at $p_{0}$ with the cuclidean $(n-\mathbf{1})$-plane $\mathbf{E}_{0}$ in $\mathbf{R}^{n}$ and via stereographic projection from $p_{\infty}$ onto $\mathrm{E}_{0}$, we identify $\mathrm{S}^{n-1}$ with the Moebius space $\mathrm{E}_{0} \cup \infty$. By abuse of notation, we regard $\psi$ as a mapping of $\mathrm{E}_{0} \cup \infty$ to itself, and similarly for the action of the Moebius group of $\mathrm{S}^{n-1}$.

Consider now the real-valued function $f$ on $\mathbf{G}_{0}$ defined by

$$
f(g)=\mathscr{C}\left(\dot{\psi}_{p_{0}}(g \mathrm{D})\right)
$$

where D is a shell in $\mathrm{E}_{0}$ and $\mathscr{C}$ denotes the conformal capacity. We regard $\mathrm{E}_{0}$ as a vector space with origin at $p_{0}$. By definition of the differential, $\psi_{p_{0}}(x)=\lim _{t \rightarrow 0} t^{-1} \psi(t x)$ for all $x \in \mathrm{E}_{0}$ and the convergence is uniform on compact subsets of $\mathrm{E}_{0}$. It follows from the last remark of Section 6, that the shell $\psi_{p}(g \mathrm{D})$ is the limit in the space of shells $\lim _{t \rightarrow 0} t^{-1} \psi(\operatorname{tg} \mathrm{D})$. By Theorem (6.I), conformal capacity depends continuously on shells, so that

$$
f(g)=\lim _{t \rightarrow 0} \mathscr{C}\left(t^{-1} \psi(t g \mathrm{D})\right)
$$

The homotheties $a_{i}: x \rightarrow t x$, where $x \in \mathrm{E}_{0}$ and $0<t \leqslant \mathrm{I}$ make up the semi-group $\mathrm{A}^{+}\left(=\mathrm{A}_{p_{0}, p_{\infty}}^{+}\right) . \quad$ Since $\Gamma_{0} \mathrm{~A}^{+}$is topologically dense in $\mathrm{G}_{0}$, we may select a sequence of elements $\gamma_{n} \in \Gamma$ and elements $t_{n}$ such that $\gamma_{n} a_{t_{n}} \rightarrow g^{-1}$ and $t_{n} \rightarrow 0$. Set $a_{n}=a_{t_{n}}(n=\mathrm{I}, 2, \ldots)$. Then

$$
\begin{aligned}
f(g) & =\lim _{t \rightarrow \infty} \mathscr{C}\left(t^{-1} \psi(\operatorname{tg} \mathrm{D})\right) \\
& =\lim _{t \rightarrow \infty} \mathscr{C}(\psi(\operatorname{tg} \mathrm{D})) \\
& =\lim _{t \rightarrow \infty} \mathscr{C}(\theta(\gamma) \psi(\operatorname{tg} \mathrm{D}))
\end{aligned}
$$

for any $\gamma \in \Gamma$, since conformal capacity is invariant under Moebius transformations. We have $\theta(\gamma) \circ \psi=\psi \circ \gamma$ for all $\gamma \in \Gamma$ and therefore we may write

$$
\begin{aligned}
f(g) & =\lim _{n \rightarrow \infty} \mathscr{C}\left(\psi\left(\gamma_{n} t_{n} g \mathrm{D}\right)\right) \\
& =\lim _{n \rightarrow \infty} \mathscr{C}\left(\psi\left(\gamma_{n} a_{n} g \mathrm{D}\right)\right) \\
& =\mathscr{C}(\psi(\mathrm{D})) .
\end{aligned}
$$

In particular

$$
\mathscr{C}\left(\dot{\psi}_{p_{\mathrm{o}}}(k \mathrm{D})\right)=\mathscr{C}\left(\dot{\psi}_{p_{0}}(\mathrm{D})\right)
$$

for all elements $k$ in the rotation subgroup $\mathrm{SO}(n-\mathrm{I}, \mathbf{R})$ of the group $\mathrm{G}_{0}$ which keeps fixed the point $p_{0}$.

Lemma (12.1). - Let T be a linear transformation of $\mathbf{R}^{n-1}$ such $\mathscr{C}(\mathrm{T}(k \mathrm{D}))=\mathscr{C}(\mathrm{T}(\mathrm{D}))$ for every $k \in \mathrm{SO}(n-\mathrm{I}, \mathbf{R})$ and for every shell D in $\mathbf{R}^{n-1}$. Then T is conformal.

Proof. - We can write $\mathrm{T}=\mathrm{T}_{1} \mathrm{~T}_{2}$ where $\mathrm{T}_{1}$ is orthogonal and $\mathrm{T}_{2}$ is a positive definite self-adjoint automorphism of $\mathbf{R}^{n-1}$. Since $\mathrm{T}_{1}$ is conformal, it suffices to prove the lemma for $T=T_{2}$. Introducing an orthonormal base of eigenvectors for $T_{2}$, we can assume that in fact T is the transformation

$$
\left(x_{1}, \ldots, x_{n-1}\right) \rightarrow\left(\lambda_{1} x_{1}, \ldots, \lambda_{n-1} x_{n-1}\right)
$$

with $\lambda_{1}=\sup \left\{\lambda_{1}, \ldots, \lambda_{n-1}\right\}$ and $\lambda_{2}=\inf \left\{\lambda_{1}, \ldots, \lambda_{n-1}\right\}$ if $n-1 \geqslant 2$. Composing $T$ with $\lambda_{1}^{-1}$ times the identity transformation, we can assume moreover that $\lambda_{1}=1$. It
remains only to show, in case $n-1 \geqslant 2$, that $\lambda_{2}=1$. We set $\lambda=\lambda_{2}$ and we know $\lambda \leqslant \mathrm{I}$.
Consider now the shell $\mathrm{D}(a)$ whose complement consists of the line segment $0 \leqslant x_{2} \leqslant a$, $0=x_{1}=x_{3}=\ldots=x_{n}$, and the infinite line segment $\mathrm{I} \leqslant x_{1}<\infty, 0=x_{2}=\ldots=x_{n}$. Let $k$ be the rotation by $90^{\circ}$ in the $x_{1}, x_{2}$ plane which leaves fixed the $x_{3}, \ldots, x_{n}$ axes. Then for all $0 \leqslant a<\infty, \mathrm{T}(k \mathrm{D}(a))$ is similar to $\mathrm{D}\left(\lambda^{-1} a\right)$, whereas $\mathrm{T}(\mathrm{D}(a))$ is similar to $\mathrm{D}(\lambda a)$. By hypothesis, $\mathscr{C}\left(\mathrm{D}\left(\lambda^{-1} a\right)\right)=\mathscr{C}(\mathrm{D}(\lambda a))$, that is $\mathscr{C}(\mathrm{D}(a))=\mathscr{C}\left(\mathrm{D}\left(\lambda^{2} a\right)\right)$ for all $a>0$. Hence $\mathscr{C}(\mathrm{D}(\mathrm{I}))=\mathscr{C}\left(\mathrm{D}\left(\lambda^{2 n}\right)\right)$. If $\lambda<\mathrm{I}$, we find $\mathscr{C}(\mathrm{D}(\mathrm{I}))=\lim _{n \rightarrow \infty} \mathscr{C}\left(\mathrm{D}\left(\lambda^{2 n}\right)\right)=\mathscr{C}(\mathrm{D}(\mathrm{o}))=0$, since $\mathrm{D}(\mathrm{o})$ has a point as one of its complementary components, by Theorem (6.1) and the remark following Lemma (7.1). But $\mathscr{C}(\mathrm{D}(\mathrm{I})) \neq 0$ by Lemma (4.2). This contradiction shows that $\lambda=1$. Hence $T$ is conformal.

It follows from Lemma (I2.I) that $\dot{\psi}_{p_{0}}$ is conformal and thus $\dot{\psi}_{p}$ is conformal for almost all $p \in \mathrm{~S}^{n-1}$. The quasi-conformal mapping $\psi$ is therefore I -quasi-conformal in the sense of Section 9 .
6. Lemma (12.2). - Let $\psi: \mathrm{S}^{n-1} \rightarrow \mathrm{~S}^{n-1}$ be a $\mathbf{1}$-quasi-conformal mapping. Then $\psi$ is a Moebius transformation if $n-\mathrm{I} \geqslant 2$.

Proof. - Composing $\psi$ with a Moebius transformation, we may assume that for some point $\mathrm{S}^{n-1}$ denoted $p_{\infty}$, we have $\psi\left(p_{\infty}\right)=p_{\infty}, \psi$ is differentiable at $p_{\infty}$, and the differential of $\psi$ at $p_{\infty}$ is the identity transformation. We shall prove after this normalization that $\psi$ is an isometry of $\mathrm{S}^{n-1}$.

Let $\pi: S^{n-1} \rightarrow \mathbf{R}^{n-1} \cup \infty$ denote stereographic projection from $p_{\infty}$. Set

$$
\xi=\pi \psi \pi^{-1} .
$$

Then $\xi$ is a I-quasi-conformal mapping of Moebius space with $\xi(\infty)=\infty$. For any point $p \in \mathbf{R}^{n-1}$ and positive real number $a$, set

$$
\begin{aligned}
l(p, a) & =\inf \{|\xi(q)-\xi(p)| ;|q-p|=a\} \\
\mathrm{L}(p, a) & =\sup \{|\xi(q)-\xi(p)| ;|q-p|=a\} \\
\mathrm{H}(p) & =\varlimsup_{a \rightarrow 0} \mathrm{~L}(p, a) / l(p, a) \\
\mathrm{I}(p) & =\varlimsup_{a \rightarrow 0} \mathrm{~L}(p, a) / a .
\end{aligned}
$$

For the point $\infty$, we define for any positive $b$

$$
\begin{aligned}
l(\infty, b) & =\sup \{|\xi(q)| ;|q|=b\} \\
\mathbf{L}(\infty, b) & =\inf \{|\xi(q)| ;|q|=b\}
\end{aligned}
$$

Since $\xi$ is I-quasi-conformal, we know that $H(p)=\mathrm{I}$ almost everywhere. The fact that $\dot{\psi}_{p_{\infty}}$ is the identity implies that

$$
\begin{equation*}
\lim _{b \rightarrow \infty} l(\infty, b) / b=1=\lim _{b \rightarrow \infty} \mathrm{~L}(\infty, b) / b \tag{12.1}
\end{equation*}
$$

Set $\mathbf{B}_{r}(p)=\{q ;|q-p| \leqslant r\}$ for any point $p \in \mathbf{R}^{n-1}$ and $\mathrm{B}_{r}(\infty)=\{q ;|q| \geqslant r\}$ for any positive real number $r$.

Let $p$ be any point of $\mathbf{R}^{n-1}$ at which $\mathrm{H}(p)=\mathrm{I}$. Let $\mathrm{D}_{a, b}(p)$ denote the spherical shell $\{q ; a<|q-p|<b\}$. We have

$$
\mathrm{B}_{l(p, a)}(\xi(p)) \subset \xi\left(\mathrm{B}_{a}(p)\right) \subset \mathrm{B}_{\mathrm{L}(p, a)}(\xi(p))
$$

and

$$
\mathrm{B}_{\mathrm{L}(\infty, b)}(\infty) \subset \xi\left(\mathrm{B}_{b}(\infty)\right) \subset \mathrm{B}_{l(\infty, b)}(\infty)
$$

Since $\mathrm{B}_{r}(p)$ and $\mathrm{B}_{r}(\infty)$ are the complementary components of $\mathrm{D}_{a, b}(p)$, we have

$$
\mathrm{D}_{\mathrm{L}(p, a), \mathrm{L}(\infty, b)}(\xi(p)) \subset \xi\left(\mathrm{D}_{a, b}(p)\right) \subset \mathrm{D}_{l(p, a), \ell \infty, b)}(\xi(p))
$$

and therefore by Lemmas (7.2) and (7.1), since $n-1>1$

$$
\text { (12.2) } \quad \log \mathrm{L}(\infty, b) / \mathrm{L}(p, a) \leqslant \log b / a \leqslant \log l(\infty, b) / l(p, a)
$$

Hence

$$
\frac{\mathrm{L}(\infty, b) / b}{\mathrm{~L}(p, a) / a} \leqslant \mathrm{I} \leqslant \frac{l(\infty, b) / b}{l(p, a) / a}
$$

Now let $a \rightarrow 0$ and $b \rightarrow \infty$. We get by (12.1) and the fact that $\mathrm{H}(p)=\mathrm{I}$,

$$
\frac{\mathrm{I}}{\mathrm{I}(p)} \leqslant \mathrm{I} \leqslant \frac{\mathrm{I}}{\mathrm{I}(p)}
$$

that is, $\mathrm{I}(p)=\mathrm{I}$ for almost all points $p \in \mathbf{R}^{n-1}$. The mapping $\xi$ thus preserves the length of all lines on which it is absolutely continuous. Since $\xi$ is absolutely continuous on almost all lines parallel to the coordinate axes, it is clearly an isometry. It follows at once that $\psi$ is a Moebius transformation.
7. We may regard $\psi$ as an element of the Moebius group $G=G M(n-1)$. Since

$$
\psi(\gamma p)=\theta(\gamma) \psi(p)
$$

for all $\gamma \in \Gamma$ and $p \in \mathrm{~S}^{n-1}$, we have

$$
\psi \gamma=\theta(\gamma) \psi
$$

and thus

$$
\theta(\gamma)=\psi \gamma \psi^{-1}
$$

with $\psi \in \mathrm{G}$. The proof of Theorem (12.I) is complete.
Remark. - The hypothesis $n>2$ was employed only twice in our proof, namely in the assertion that the modulus of the spherical shell $\mathrm{D}_{a, b}: a<|x|<b$ is $\log b / a$ and in Theorem (9.4). In dimension I the conformal capacity of such a domain has the constant value 2. And of course, Theorem (12.I) is false in case $n=2$. A counter example is provided by two compact Riemann surfaces which have the same topological genus but not the same conformal type. In fact, with that example in mind, we state the following corollaries of Theorem (12.1).

Corollary (12.2). -- Let Y and $\mathrm{Y}^{\prime}$ be complete Riemannian manifolds having constant negative curvature and finite volume. Assume that there is a quasi-conformal homeomorphism of Y onto $\mathrm{Y}^{\prime}$. Then there is a unique conformal mapping of Y onto $\mathrm{Y}^{\prime}$ inducing the same isomorphism
of fundamental groups, provided the dimension is greater than two. The conformal mapping may be taken as an isometry if Y and $\mathrm{Y}^{\prime}$ have the same curvature.

Proof. - Let X and $\mathrm{X}^{\prime}$ denote the simply connected covering manifolds of Y and $\mathrm{Y}^{\prime}$ respectively. Then, as is well-known, X and $\mathrm{X}^{\prime}$ are hyperbolic $n$-spaces and up to isometries are completely characterized by a single invariant, namely the curvature. By modifying the metric via a multiplicative constant at all points, the curvatures can be made equal. In other words, X and $\mathrm{X}^{\prime}$ are conformally equivalent. We may assume without loss of generality that $\mathrm{X}=\mathrm{X}^{\prime}$. Let G denote the group of isometries of $X$, and let $\Gamma$ and $\Gamma^{\prime}$ denote the fundamental groups of $Y$ and $Y^{\prime}$ respectively. A quasiconformal mapping of $Y$ onto $\mathrm{Y}^{\prime}$ induces an isomorphism $\theta: \Gamma \rightarrow \Gamma^{\prime}$ and a quasiconformal mapping $\varphi: X \rightarrow X$ such that

$$
\varphi(\gamma x)=\theta(\gamma) \varphi(x)
$$

for all $\gamma \in \Gamma, x \in \mathrm{X}$. The result now follows from Theorem (12.I). Namely, if

$$
\psi \gamma \psi^{-1}=\theta(\gamma), \quad \psi \in \mathrm{G}
$$

then $\psi(\gamma x)=\theta(\gamma) \psi(x)$ for all $\gamma \in \Gamma, x \in \mathrm{X}$ and $\psi$ induces a conformal mapping of Y to $\mathrm{Y}^{\prime}$ - in fact an isometry.

To prove that $\psi$ is unique, suppose that $\psi$ and $\psi^{\prime}$ are elements of $\mathbf{G}$ satisfying

$$
\psi \gamma=\theta(\gamma) \psi \quad \text { and } \quad \psi^{\prime} \gamma=\theta(\gamma) \psi^{\prime}
$$

for all $\gamma \in \Gamma$. Set $\zeta=\psi^{-1} \psi^{\prime}$. Then

$$
\zeta \gamma=\gamma \zeta
$$

for all $\gamma \in \Gamma$. Since $\Gamma \backslash G$ has finite measure, it follows that $\zeta$ is central in $G$ (cf. A. Borel [2], or G. D. Mostow [14]). Since the Moebius group has trivial center, $\zeta=1$ and $\psi=\psi^{\prime}$.

Corollary (12.3). - Any two compact Riemannian manifolds of the same constant negative curvature are isometric if they are diffeomorphic, provided the dimension is greater than two.

Proof. - This follows immediately from Corollary (12.I) once we remark that a diffeomorphism of compact Riemannian manifolds is quasi-conformal.

Remark. - Let $\varphi_{0}$ and $\varphi_{1}$ be continuous mappings of hyperbolic $n$-space X into itself such that for all $x \in X$ and $\gamma \in \Gamma$

$$
\begin{aligned}
& \varphi_{0}(\gamma x)=\theta(\gamma) \varphi_{0}(x) \\
& \varphi_{1}(\gamma x)=\theta(\gamma) \varphi_{1}(x)
\end{aligned}
$$

It is not hard to see that there is a unique geodesic joining any two points in X . Let $\varphi_{t}(x)$ denote the point on the geodesic segment $\left[\varphi_{0}(x), \varphi_{1}(x)\right]$ joining $\varphi_{0}(x)$ to $\varphi_{1}(x)$ that divides it in the ratio $t /(\mathrm{I}-t)$ where $0 \leqslant t \leqslant 1$. Then $\varphi_{t}: x \rightarrow \varphi_{t}(x)$ is continuous and we have for all $0 \leqslant t \leqslant \mathrm{I}, \gamma \in \Gamma, x \in \mathrm{X}$

$$
\begin{equation*}
\varphi_{t}(\gamma x)=\theta(\gamma) \varphi_{t}(x) \tag{*}
\end{equation*}
$$

since $\theta(\gamma)$ is an isometry of X sending $\left[\varphi_{0}(x), \varphi_{1}(x)\right]$ to $\left[\varphi_{0}(\gamma x), \varphi_{1}(\gamma x)\right]$. We apply this observation to two continuous mapping $\Phi_{0}: \Gamma \backslash \mathrm{X} \rightarrow \Gamma^{\prime} \backslash \mathrm{X}$ and $\Phi_{1}: \Gamma \backslash \mathrm{X} \rightarrow \Gamma^{\prime} \backslash \mathrm{X}$ which induce the same homomorphism $\theta: \Gamma \rightarrow \Gamma^{\prime}$ of the fundamental groups. We obtain from condition (*) a deformation $\Phi_{i}: \Gamma \backslash \mathrm{X} \rightarrow \Gamma^{\prime} \backslash \mathrm{X}(0 \leqslant t \leqslant 1)$ of $\Phi_{0}$ to $\Phi_{1}$. Thus the homotopy class of a continuous mapping of hyperbolic space forms is uniquely determined by the induced homomorphism of the fundamental groups. In particular, the conformal mapping $\psi$ of Corollary (12.2) is homotopic to the quasi-conformal mapping $\varphi$.

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