## SUMMABILITY METHODS ON MATRIX SPACES

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§1. Introduction. The matrix spaces under consideration are the four main types of irreducible bounded symmetric domains given by Cartan (5). Let $z=\left(z_{j k}\right)$ be a matrix of complex numbers, $z^{\prime}$ its transpose, $z^{*}$ its conjugate transpose and $I=I^{(n)}$ the identity matrix of order $n$. Then the first three types are defined by

$$
\begin{equation*}
D=\left[z \mid I-z z^{*}>0\right] \tag{1}
\end{equation*}
$$

where $z$ is an $n$ by matrix $(n \leqslant m$ ), a symmetric or a skew-symmetric matrix of order $n$ (16). The fourth type is the set of complex spheres satisfying

$$
\left|z^{\prime} z\right|<1,1-2 z^{*} z+\left|z^{\prime} z\right|^{2}>0
$$

where $z$ is an $n$ by 1 matrix. It is known that each of these domains possesses a distinguished boundary $B$ which in the first three cases is given by

$$
\begin{equation*}
B=\left[u \mid u u^{*}=I\right] . \tag{2}
\end{equation*}
$$

(In the case of skew symmetric matrices the distinguished boundary is given by (2) only if $n$ is even.)

In §2 we consider the following problem for the first type of domain with $m=n$, in which case $u$ is a unitary matrix, the (real) dimension of $B$ is $n^{2}$ and of $D$ is $2 n^{2}$. Let $f(u)$ be a real integrable function defined on $B$ and consider the integral operator

$$
\begin{equation*}
I(f, z)=\int_{B} P(z, u) f(u) d V, \tag{3}
\end{equation*}
$$

where $P(z, u)$ is the Poisson kernel (14)

$$
\begin{equation*}
P(z, u)=V^{-1} \operatorname{det}^{n}\left(I-z u^{*}\right)^{-1}\left(I-z z^{*}\right)\left(I-u z^{*}\right)^{-1} \tag{4}
\end{equation*}
$$

$V$ is the Euclidean volume of $B$, and $d V$ the Euclidean volume element. It is known that $I(f, z)$ is a harmonic function of $z$ if $I-z z^{*}>0$ or $I-z z^{*}<0$. (I proved this fact in (14) for $z \in D$ but the proof is valid for all $z$ and $u \in B$ for which $\operatorname{det}\left(I-z u^{*}\right) \neq 0$. It is easily proved that $\operatorname{det}\left(I-z u^{*}\right) \neq 0$ for $u \in B$ and all $z$ such that $I-z z^{*}>0$ or $I-z z^{*}<0$.) Here a harmonic function is a function of class $C^{2}$, which satisfies on $D$ the Laplace equation corresponding to the invariant metric of $D$, that is, the metric invariant with respect to the group of 1 to 1 analytic transformations mapping $D$ onto itself (14). This invariant metric is given by

$$
d s^{2}=\sigma\left[\left(I-z z^{*}\right)^{-1} d z\left(I-z^{*} z\right)^{-1} d z^{*}\right],
$$

[^0]where $\sigma(A)$ is the trace of the matrix $A$ and $d z=\left(d z_{j k}\right)$, and the corresponding Laplace equation is
$$
4 \sigma\left[\bar{\partial}\left(I-z^{*} z\right) \partial^{\prime}\left(I-z z^{*}\right)\right]=0, \quad \partial=\left(\partial / \partial z_{j k}\right)
$$

It has been proved by Hua and Lowdenslager that given a real function $f$, continuous on $B$, there exists a function $F$, harmonic on $D$, such that $F(z) \rightarrow$ $f\left(u_{0}\right)$ as $z \rightarrow u_{0} \in B$ radially, that is, along the set $\rho u_{0}, 0 \leqslant \rho<1(6 ; 9)$. Further, if $F$ is continuous on the closure $\bar{D}$ of $D$ and satisfies certain other conditions due to Lowdenslager on the boundary of $D$ other than $B$, then $F$ is unique (8). Now for the particular case of the unit circle, $z \bar{z}=1$, if we merely assume that $f$ is integrable on it, then $I\left(f, \rho u_{0}\right) \rightarrow f\left(u_{0}\right)$ as $\rho \rightarrow 1(0 \leqslant \rho<1)$ (20); this method of approach is known as Abel-Poisson summability of Fourier series. We prove this result for matrix spaces. (See note added in proof).

In §3 we consider for the first type of domain $(n \leqslant m)$ some properties of complete orthonormal systems (CONS) of complex homogeneous polynomials defined on $D$. The space $D$ is circular with center at $z=0$, that is, if $z \in D$, then $e^{i \theta} z \in D$ for $0 \leqslant \theta<2 \pi$. Hence any two powers

$$
\begin{equation*}
P(z)=\prod_{j, k} z_{j k}^{s_{j k}}, \quad Q(z)=\prod_{j, k} z_{j k}^{t_{j k}}, \tag{5}
\end{equation*}
$$

$s_{j k}, t_{j k}$ non-negative integers, for which

$$
\begin{equation*}
\sum_{j, k} s_{j k} \neq \sum_{j, k} t_{j k}, \tag{6}
\end{equation*}
$$

are orthogonal, that is,

$$
\begin{equation*}
(P, Q)=\int_{D} P(z) \bar{Q}(z) d W=0 \tag{7}
\end{equation*}
$$

( $d W$ is Euclidean volume element on $D$ ) (6). Also if $f \in$ class $L^{2}$ on $D$, which means that $f$ is single-valued and analytic on $D$ and has finite norm $\|f\|=$ $[(f, f)]^{\frac{1}{2}}$ (2), then the set $\{P\}$ is complete with respect to functions of class $L^{2}$ (6).

Here we refine conditions (6) to show that $(P, Q) \neq 0$ implies
(8) $\quad \sum_{k} s_{\nu k}=\sum_{k} t_{\nu k}, \quad \sum_{j} s_{j \mu}=\sum_{j} t_{j \mu} \quad(\nu=1, \ldots, n ; \mu=1, \ldots, m)$.

By means of (8) the set of powers $\{P\}$ is subdivided into disjoint subsets whose members need not be orthogonal to each other. The elements of a subset are made into an orthonormal set by the Gram Schmidt formulas, thus giving a CONS of homogeneous polynomials $\{\phi\}$ on $D$. We note that Hua has constructed a CONS of functions of class $L^{2}$ on $D$ using representation theory (6).

In $\S 4$ applications of the CONS $\left\{\phi_{\nu}\right\}$ are given. First an Abelian theorem is obtained and then a Tauberian theorem for the orthogonal series $\sum a_{\nu} \phi_{\nu}$ as $z$ approaches $[I, 0]$ of $B$ along the matrix $[r, 0]$ where $r$ is the diagonal matrix $\left[r_{1}, \ldots, r_{n}\right], 0 \leqslant r_{j}<1$. Next a Cauchy's inequality is obtained for the Fourier coefficients $a_{\nu}$. Finally two mean value theorems, which generalize analogous theorems for the unit circle, are proved.

## §2. Poisson summability.

1. Reduction of integral (1.3) to.normal form. Rauch outlined this reduction to me. The transformation

$$
\begin{equation*}
w=z u_{0}^{-1}, \quad u_{0} u_{0}^{*}=I, \tag{1}
\end{equation*}
$$

takes $z=u_{0}$ into $w=I$ and also leaves $D$ and $B$ invariant since under it $I-w w^{*}=I-z z^{*}$. Also if $u \rightarrow v$ under (1) $P(z, u) \rightarrow P(w, v)$ and

$$
d V_{u}=\frac{\dot{u}}{\operatorname{det}^{n} u}=\frac{\partial(v)}{\partial(u)} \frac{\dot{v}}{\operatorname{det}^{n} v \operatorname{det}^{n} u_{0}},
$$

where

$$
\dot{u}=(-1)^{-\frac{1}{4} n(n+1)} \prod_{j, k} d u_{j k}
$$

(14). Now $d u=d v u_{0}$, the Jacobian of which is $\partial(u) / \partial(v)=\operatorname{det}^{n} u_{0}$ (3). Thus $d V_{u} \rightarrow d V_{v}$. Also $f(u) \rightarrow f\left(v u_{0}\right)=f_{0}(v)$ so that $I(f, z) \rightarrow I\left(f_{0}, w\right)$.

If $w \rightarrow I$ along the set of points $\rho I(0 \leqslant \rho<1)$, then

$$
\operatorname{det}\left(I-w w^{*}\right)=\left(1-\rho^{2}\right)^{n}
$$

and

$$
\begin{aligned}
Q & =\left(I-w v^{*}\right)\left(I-v w^{*}\right)=I-w v^{*}-v w^{*}+w w^{*} \\
& =I\left(1+\rho^{2}\right)-\rho\left(v+v^{*}\right) .
\end{aligned}
$$

Now $v$ is unitary equivalent to a diagonal matrix $v_{D}$ which is also unitary ( $\mathbf{1 0}$, Theorem 41.41), that is,

$$
\begin{equation*}
v=U^{*} v_{D} U \tag{2}
\end{equation*}
$$

Thus if $v_{D}=\left[d_{1}, \ldots, d_{n}\right]$, then $d_{j} \bar{d}_{j}=1$ and we can write $d_{j}=e^{i \theta_{j}}\left(0 \leqslant \theta_{j}<\right.$ $2 \pi)$. Hence

$$
Q=U^{*}\left[I\left(1+\rho^{2}\right)-\rho\left(v_{D}^{*}+v_{D}\right)\right] U
$$

and

$$
\operatorname{det} Q=\operatorname{det}\left[I\left(1+\rho^{2}\right)-\rho\left(v_{D}^{*}+v_{D}\right)\right]=\prod_{j=1}^{n}\left(1-2 \rho \cos \theta_{j}+\rho^{2}\right) .
$$

Let

$$
\begin{equation*}
p\left(\rho, \theta_{j}\right)=\frac{1-\rho^{2}}{1-2 \rho \cos \theta_{j}+\rho^{2}}, \tag{3}
\end{equation*}
$$

which is $2 \pi$ times the Poisson kernel for the unit circle. Then

$$
\begin{equation*}
P(\rho I, v)=V^{-1} \prod_{j=1}^{n} p^{n}\left(\rho, \theta_{j}\right) \tag{4}
\end{equation*}
$$

and

$$
\begin{equation*}
I\left(f_{0}, \rho I\right)=V^{-1} \int_{B} \prod_{j} p^{n}\left(\rho, \theta_{j}\right) f_{0}(v) d V_{v} . \tag{5}
\end{equation*}
$$

According to Weyl (19, p. 197) if (2) holds, then

$$
\begin{equation*}
d V_{v}=\left[d V_{U}\right] \Delta \bar{\Delta} d \theta_{1} \ldots d \theta_{n} \tag{6}
\end{equation*}
$$

where

$$
\begin{equation*}
\Delta \bar{\Delta}=\prod_{j<k}\left|e^{i \theta_{j}}-e^{i \theta_{k}}\right|^{2}=4 \prod_{j<k} \sin ^{2} \frac{1}{2}\left(\theta_{j}-\theta_{k}\right) \tag{7}
\end{equation*}
$$

and $d V_{0}=\left[d V_{U}\right]$ is a constant times the Euclidean volume element on ${ }^{\text {r }}$ the other $n^{2}-n$ parameters defining $B$. Let $B_{0}$ be this part of $B$. Now $P(\rho I, v)$, considered as a function of $v$, is independent of the other $n^{2}-n$ parameters and hence

$$
\begin{equation*}
I\left(f_{0}, \rho I\right)=\int_{0}^{2 \pi} \cdots \int_{0}^{2 \pi} P(\rho I, v) F(\theta) \Delta \bar{\Delta} d \theta_{1} \ldots d \theta_{n} \tag{8}
\end{equation*}
$$

where

$$
\begin{equation*}
F(\theta)=F\left(\theta_{1}, \ldots, \theta_{n}\right)=\int_{B_{0}} f_{0}(v) d V_{0} . \tag{9}
\end{equation*}
$$

2. Convergence theorem for (8). It is sufficient to consider (8) for $n=2$, in which case replacing $\theta_{1}$ by $s$ and $\theta_{2}$ by $t$

$$
\begin{equation*}
I\left(f_{0}, \rho I\right)=4 V^{-1} \int_{0}^{2 \pi} \int_{0}^{2 \pi} p^{2}(\rho, s) p^{2}(\rho, t) \sin ^{2} \frac{1}{2}(s-t) F(s, t) d s d t \tag{10}
\end{equation*}
$$

and we consider $\lim _{\rho \rightarrow 1} I\left(f_{0}, \rho I\right)$. Subtracting $S$ from each side we reduce (10) by well-known methods in Fourier series (20) to a consideration of

$$
\begin{equation*}
I=4 V^{-1} \int_{0}^{\pi} \int_{0}^{\pi} p^{2}(\rho, s) p^{2}(\rho, t) \sin ^{2} \frac{1}{2}(s-t) G(s, t) d s d t \tag{11}
\end{equation*}
$$

and

$$
\begin{equation*}
I_{1}=4 V^{-1} \int_{0}^{\pi} \int_{0}^{\pi} p^{2}(\rho, s) p^{2}(\rho, t) \sin ^{2} \frac{1}{2}(s+t) G(2 \pi-s, t) d s d t \tag{11a}
\end{equation*}
$$

where

$$
\begin{equation*}
G(s, t)=F(s, t)+F(2 \pi-s, 2 \pi-t)-2 V_{0} S, \tag{12}
\end{equation*}
$$

$V_{0}$ being the volume of $B_{0}$. We show that under certain hypotheses on $G(s, t)$, $I \rightarrow 0$ as $\rho \rightarrow 1$ and the proof is similar for $I_{1}$. In (11) integrate by parts with respect to $s$ and $t$. Since

$$
H(s, t)=\iint G(s, t) d s d t
$$

is zero for $s=0$ or $t=0, \sin ^{2} \frac{1}{2}(s-t)=0$ for $s=t=\pi, \sin ^{2} \frac{1}{2}(\pi-\theta)=$ $\cos ^{2} \frac{1}{2} \theta$ and $p(\rho, \pi)=(1-\rho) /(1+\rho)$, we get

$$
\begin{align*}
I= & -\left(\frac{1-\rho}{1+\rho}\right)^{2} \int_{0}^{\pi} \frac{\partial}{\partial t}\left[p^{2}(\rho, t) \cos ^{2} \frac{1}{2} t\right] H(\pi, t) d t  \tag{13}\\
& -\left(\frac{1-\rho}{1+\rho}\right)^{2} \int_{0}^{\pi} \frac{\partial}{\partial s}\left[p^{2}(\rho, s) \cos ^{2} \frac{1}{2} s\right] H(s, \pi) d s \\
& +\int_{0}^{\pi} \int_{0}^{\pi} \frac{\partial^{2}}{\partial s \partial t}\left[p^{2}(\rho, s) p^{2}(\rho, t) \sin ^{2} \frac{1}{2}(s-t)\right] H(s, t) d s d t \\
\equiv & I_{1}+I_{2}+I_{3} .
\end{align*}
$$

The following inequalities are used in considering $I_{1}, I_{2}$, and $I_{3}$ (20):
(i)

$$
\begin{equation*}
0 \leqslant p(\rho, \theta) \leqslant \frac{1+\rho}{1-\rho} \tag{14}
\end{equation*}
$$

(ii) $\quad 0 \leqslant p(\rho, \theta) \leqslant \frac{1-\rho^{2}}{4 \sin ^{2} \frac{1}{2} \theta}$
(iii) $\lim _{\theta \rightarrow 0} \frac{\theta}{\sin \theta}=1$
(iv) $\quad \int_{0}^{\pi} p(\rho, \theta) d \theta=\pi$

$$
\begin{equation*}
\int_{\delta}^{\pi} p(\rho, \theta) d \theta=o(1) \quad \text { as } \quad \rho \rightarrow 1 \quad \text { for } \quad 0<\delta \leqslant \pi \tag{v}
\end{equation*}
$$

Also

$$
\frac{\partial}{\partial \theta} p(\rho, \theta)=2 \rho\left(1-\rho^{2}\right) \sin \theta d^{-2}(\rho, \theta)
$$

where

$$
d(\rho, \theta)=1-2 \rho \cos \theta+\rho^{2}=\frac{1-\rho^{2}}{p(\rho, \theta)}
$$

and

$$
\begin{aligned}
& \frac{\partial^{2}}{\partial s \partial t}\left[p^{2}(\rho, s) p^{2}(\rho, t) \sin ^{2} \frac{1}{2}(s-t)\right]=p(\rho, s) p(\rho, t) . \\
& \quad\left[16 \rho^{2}\left(1-\rho^{2}\right)^{2} \sin s \sin t \sin ^{2} \frac{1}{2}(s-t) d^{-2}(\rho, s) d^{-2}(\rho, t)-\frac{1}{2} p(\rho, s) p(\rho, t)\right. \\
& \quad . \cos (s-t)-2 \rho\left(1-\rho^{2}\right) \sin s \sin (s-t) p(\rho, t) d^{-2}(\rho, s) \\
& \left.\quad+2 \rho\left(1-\rho^{2}\right) \sin t \sin (s-t) p(\rho, s) d^{-2}(\rho, t)\right] .
\end{aligned}
$$

Concerning the function $G(s, t)$ we assume that

$$
\begin{align*}
& \lim _{h . k \rightarrow 0} \frac{1}{h k} \int_{0}^{h} \int_{0}^{k}|G(s, t)| d s d t=0  \tag{15}\\
& \int_{0}^{h} \int_{0}^{k}|G(s, t)||d s d t| \leqslant C|h k|, \quad(0<|h|,|k| \leqslant \pi),
\end{align*}
$$

$C$ an absolute constant. (See (11) where these conditions were used in a similar connection.) Then using (15) we find

$$
\begin{aligned}
\left|I_{1}\right| \leqslant \pi\left(1-\rho^{2}\right) C \int_{0}^{\pi}\left|2 p(\rho, t) \cos ^{2} \frac{1}{2} t \partial p(\rho, t) / \partial \mathrm{t}-p^{2}(\rho, t) \sin \frac{1}{2} t \cos \frac{1}{2} t\right| t d t
\end{aligned} \quad \begin{aligned}
& =0\left((1-\rho) \int_{0}^{\pi} p(\rho, t) d t\right)=0(1-\rho)
\end{aligned}
$$

and similarly for $I_{2}$. Now

$$
I_{3}=\int_{0}^{\delta} \int_{0}^{\delta}+\left\{\int_{0}^{\delta} \int_{\delta}^{\pi}+\int_{\delta}^{\pi} \int_{0}^{\delta}+\int_{\delta}^{\pi} \int_{\delta}^{\pi}\right\} \equiv I_{31}+I_{32}
$$

By (15) it is found that

$$
\left|I_{31}\right| \leqslant \epsilon \int_{0}^{\pi} \int_{0}^{\pi} p(\rho, s) p(\rho, t) d s d t
$$

where $\epsilon$ is an arbitrary positive number, and

$$
I_{32}=0\left(\frac{1-\rho}{\sin ^{2} \frac{\rho}{2} \delta}\right)
$$

Consequently given $\epsilon>0$ choose $\delta>0$ so that

$$
\left|\int_{0}^{s} \int_{0}^{t} G(s, t) d s d t\right|<\epsilon|s t|,
$$

if $|s|<\delta$ and $|t|<\delta$. With fixed $\delta$ choose $1-\rho$ sufficiently small. Then $I=0(\epsilon)$ for $\rho$ sufficiently close to 1 . Consequently we have proved

Theorem 2.1. Let $F(s, t)$ be an integrable function. If $G(s, t)=F(s, t)+$ $F(2 \pi-s, 2 \pi-t)-2 V_{0} S$ satisfies conditions (15), then

$$
\lim _{\rho \rightarrow 1^{-}} 4 V^{-1} \int_{0}^{2 \pi} \int_{0}^{2 \pi} p^{2}(\rho, s) p^{2}(\rho, t) \sin ^{2} \frac{1}{2}(s-t) F(s, t) d s d t=S .
$$

For $n>2$ it would be sufficient to assume that

$$
\begin{align*}
& \lim _{\theta_{j} \rightarrow 0} \frac{1}{\theta_{1} \ldots \theta_{n}} \int_{0}^{\theta_{n}} \ldots \int_{0}^{\rho_{n}}|G(\theta)| d \theta_{1} \ldots d \theta_{n}=0  \tag{16}\\
& \int_{0}^{\theta_{1}} \ldots \int_{0}^{\theta_{n}}|G(\theta)|\left|d \theta_{1} \ldots d \theta_{n}\right| \leqslant C\left|\theta_{1} \ldots \theta_{n}\right|,\left(0<\left|\theta_{j}\right| \leqslant \pi, j=1, \ldots, n\right),
\end{align*}
$$

where $G(\theta)$ is defined similarly to $G(s, t)$. We obtain
Theorem 2.2. Let $f$ be an integrable function on the unitary group $B$ such that the function $F(\theta)$ defined by (9), where $f_{0}$ is the transform of $f$ under (1), satisfies (16). Then the Poisson integral (1.3) has a limit if $z$ approaches the point $u_{0}$ on $B$ radially.

## §3. Complete orthonormal systems on D. Orthogonal developments.

1. Integration over $D$. Let $z$ be an $n$ by $m$ matrix $(n \leqslant m), z_{p}$ its $p$ th row and $Z_{p}$ the submatrix consisting of the first $p$ rows $(p=1, \ldots, n)$. The inner product ( $f, g$ ) defined by (1.7) may be transformed into an iterated integral over the product of $n$ hyperspheres by a procedure due to Hua (12), giving

$$
\begin{equation*}
(f, g)=\left(\frac{1}{2 i}\right)^{m n} \int_{w_{1} w_{1}^{*}<1}\left(1-w_{1} w_{1}^{*}\right)^{n-1} \dot{w}_{1} \ldots \int_{w_{n} w_{n}^{*}<1} f \bar{g} \dot{w}_{n} \tag{1}
\end{equation*}
$$

where

$$
\begin{align*}
& \dot{w}_{p}=\prod_{k=1}^{m} d w_{p k} d \bar{w}_{p k}  \tag{2}\\
& w_{p}=z_{p} \Gamma_{p-1} \\
& z_{p}=w_{p} \Gamma_{p-1}^{-1} \quad\left(p=1, \ldots, n, \Gamma_{0}=1\right)
\end{align*}
$$

and $\Gamma_{p-1}$ is a unique positive definite matrix such that

$$
\begin{equation*}
\Gamma_{p-1} \Gamma_{p-1}^{*}=\left(I-Z_{p-1}^{*} Z_{p-1}\right)^{-1} \tag{3}
\end{equation*}
$$

2. Construction of the matrix $\Gamma_{p-1}^{-1}(2 \leqslant p \leqslant n)$. Let $U_{p}(q)=U(q)$ be the minor formed from the first $q$ rows and columns of

$$
\begin{equation*}
U_{p}=U=\left(u_{j k}\right)=I-Z_{p-1}^{*} Z_{p-1} \tag{4}
\end{equation*}
$$

$\mathscr{U}_{p}(q)=\mathscr{U}(q)$ the corresponding submatrix and $u_{m-j}=\left(u_{m-j, 1}, \ldots, u_{m-j, m-j-1}\right)$. Since $I-Z_{p-1} Z^{*}{ }_{p-1}$ is the leading $(p-1)$ th principal submatrix of $I-z z^{*}$, it is positive definite, hence $U$ is positive definite and $U(q)$ are positive. Thus the hermitian matrix $U$ may be reduced to diagonal form by the well-known Kronecker reduction (1) whose ( $j+1$ )th step is

$$
V_{j+1} \ldots V_{1} U V_{1}^{*} \ldots V_{j+1}^{*}=\left(\begin{array}{ccc}
\mathscr{U}(m-j-1) & 0 \\
& & X_{m, m-j} \\
0 & \ddots & \\
& & \\
& & X_{m m}
\end{array}\right)
$$

$(0 \leqslant j \leqslant m-2)$, where

$$
\begin{aligned}
V_{j+1} & =\left(\begin{array}{ll}
I^{(m-j-1)} & 0 \\
-u_{m-j} \mathscr{U}^{-1}(m-j-1) I^{(j+1)} \\
0
\end{array}\right), \\
X_{m, m-j} & =U(m-j) U^{-1}(m-j-1) \\
X_{m 1} & =U(1)
\end{aligned}
$$

Also

$$
\begin{aligned}
& u_{m-j} \mathscr{U}^{-1}(m-j-1)=U^{-1}(m-j-1) \\
& \quad\left((-1)^{k+m-j+1} U\binom{1, \ldots,[k], \ldots, m-j}{1, \ldots, m-j-1}\right),(1 \leqslant k \leqslant m-j-1),
\end{aligned}
$$

where

$$
U\binom{1, \ldots,[k], \ldots, m-j}{1, \ldots, m-j-1}
$$

is the minor of $U$ formed from rows $1, \ldots, m-j$ with row $k$ omitted and columns $1, \ldots, m-j-1$.

Now $V_{j+1}^{-1}$ equals $V_{j+1}$ with the sign of the matrix $-u_{m-j} \mathscr{U}^{-1}(m-j-1)$ changed. Also $X_{m k}$ is real. Hence we can take

$$
\begin{equation*}
\Gamma_{p-1}^{*-1}=V_{1}^{-1} \ldots V_{m-1}^{-1}\left[X_{m}^{\frac{1}{2}}, \ldots, X_{m m}^{\frac{1}{2}}\right] \tag{5}
\end{equation*}
$$

By an inductive proof we may show that

$$
\begin{equation*}
V_{1}^{-1} \ldots V_{m-1}^{-1}=\left(\frac{U\binom{j}{1}}{U(1)} \frac{U\binom{1, j}{1,2}}{U(2)} \ldots \frac{U\binom{1,2, \ldots,[j-1], j}{1,2, \ldots, j-1}}{U(j-1)} 10 \ldots 0\right) \tag{6}
\end{equation*}
$$

$(1 \leqslant j \leqslant m)$. (For details of the proof cf. (15).) From (4) it follows that

$$
\begin{align*}
U\binom{1, \ldots,[i], r}{1, \ldots, i} & =U_{p}\binom{1, \ldots,[i], r}{1, \ldots, i}  \tag{7}\\
& =\operatorname{det}\left(\delta_{j k}-\sum_{l=1}^{p-1} \bar{z}_{l j} z_{l k}\right) \\
& (j=1, \ldots, i-1, r ; k=1, \ldots, i ; i+1 \leqslant r \leqslant m, 1 \leqslant i \leqslant m-1)
\end{align*}
$$

3. Formula for $z_{p r}$. From (2), (5), and (6) follows

$$
\begin{align*}
& z_{p r}=\sum_{i=1}^{r-1} c_{p i} \bar{U}_{p}\binom{1, \ldots,[i], r}{1, \ldots, i} w_{p i}+c_{p r r} w_{p r}(p \geqslant 2)  \tag{8}\\
& z_{1 r}=w_{1 r}\left(1 \leqslant r \leqslant m ; \sum_{i=1}^{0}=0\right),
\end{align*}
$$

where

$$
\begin{align*}
c_{p i} & =\left[U_{p}(i) U_{p}(i-1)\right]^{-\frac{1}{2}}(1 \leqslant i \leqslant r-1)  \tag{9}\\
c_{p r r} & =\left[U_{p}(r) / U_{p}(r-1)\right]^{\frac{1}{2}} \\
c_{1 i} & =c_{1 r r}=1 .
\end{align*}
$$

The formula

$$
\begin{equation*}
U_{p}(q)=\prod_{k=1}^{p-1}\left(1-\sum_{j=1}^{q} w_{k j} \bar{w}_{k j}\right) \quad(p \geqslant 2) \tag{10}
\end{equation*}
$$

holds.
We prove (10) by induction on $p$. Since by (8)

$$
U_{2}(q)=\operatorname{det}\left(I-z_{1}^{*} z_{1}\right)=\operatorname{det}\left(I-z_{1} z_{1}^{*}\right)=1-\sum_{j=1}^{q} w_{1 j} \bar{w}_{1 j},
$$

where $z_{1}=\left(z_{11}, \ldots, z_{1 q}\right),(10)$ is true for $p=2$. Assume (10) holds for $U_{p}(q)$ and prove for $U_{p+1}(q)$. Since $U_{p}(q) \neq 0$, there exists a unimodular matrix $A$ such that the matrix $\mathscr{V}_{p+1}(q)=\left(\delta_{j k}-\sum_{l=1}^{q} z_{j l} \bar{z}_{k l}\right)(1 \leqslant j, k \leqslant p)$ is equal to

$$
A\left[\mathscr{V}_{p}(q), 1-z_{p} \mathscr{U}_{p}^{-1}(q) z_{p}^{*}\right] A^{*},
$$

(12), and thus

$$
U_{p+1}(q)=\operatorname{det} \mathscr{V}_{p+1}(q)=\operatorname{det} \mathscr{V}_{p}(q)\left(1-z_{p} \mathscr{U}_{p}^{1}(q) z_{p}^{*}\right) .
$$

Since det $\mathscr{V}_{p}(q)=U_{p}(q)$, using (10) we need only consider the last factor on the right, which equals

$$
E=1-\sum_{j, k=1}^{q} z_{p j} U_{k j} \bar{z}_{p k} / U_{p}(q),
$$

where $U_{k^{j}}$ is the cofactor of the element $u_{k j}$ in the matrix $\mathscr{U}_{p}(q)$. Substitute (8) for $z_{p j}$ and $\bar{z}_{p k}$. The term $w_{p r}$ occurs when $j=r$ or when $i=r$ and $j=r+1$, $\ldots, q$ and $\bar{w}_{p \lambda}$ occurs when $k=\lambda$ or $i=\lambda$ and $k=\lambda+1, \ldots, q$. Thus the coefficient, $C_{r \lambda}$, of $w_{p r} \bar{w}_{p \lambda}(1 \leqslant r, \lambda \leqslant q)$ is

$$
C_{r \lambda}=U_{p}^{-1}(q)\left[c_{p r} \sum_{j=r+1}^{q} \bar{U}_{p}\binom{1, \ldots,[r], j}{1, \ldots, r} D_{j \lambda}+c_{p r r} D_{r \lambda}\right]
$$

where

$$
D_{j \lambda}=c_{p \lambda} \sum_{k=\lambda+1}^{q} U_{k j} U_{p}\binom{1, \ldots,[\lambda], k}{1, \ldots, \lambda}+U_{\lambda j} c_{p \lambda \lambda .}
$$

By means of elementary properties of determinants it is not difficult to prove in case $\lambda \leqslant r$ that $D_{j \lambda}=0$ for $j=r, \lambda<r$ and $j=r+1, \ldots, q, \lambda \leqslant r$. Hence $C_{r \lambda}=0$ for $\lambda \neq r$. Also $C_{r r}=1$. A similar proof holds for $r<\lambda$. Thus

$$
E=1-\sum_{j=1}^{q} w_{p j} \bar{w}_{p j}
$$

and $U_{p+1}(q)$ has the desired form, which proves (10).
4. Structure of the CONS on D.

Theorem 3.1. ( $P, Q$ ) $\neq 0$ implies equations (1.8).
Proof. We first show that

$$
\begin{equation*}
z_{p r}=w_{p r} \sum B_{i} w_{p i} \bar{w}_{p r} \Pi\left(w_{\lambda_{1} \alpha_{1}} \bar{w}_{\lambda_{1} \beta_{1}} \ldots w_{\lambda_{l \alpha} l} \bar{w}_{\lambda_{l \beta} \beta_{l}} w_{\lambda_{i} r} \bar{w}_{\lambda_{i} \beta_{i}}\right), \tag{11}
\end{equation*}
$$

where $B_{i}$ is a function of $w_{j k} \bar{w}_{j k}(1 \leqslant j \leqslant p-1) ; \lambda_{1}, \ldots, \lambda_{l}, \lambda_{i}$ take on values in the set $1,2, \ldots, p-1 ; \alpha_{1}, \ldots, \alpha_{l}$ is a subset of $1, \ldots, i-1$ and $\left(\beta_{1}, \ldots\right.$, $\left.\beta_{l}, \beta_{i}\right)$ is a permutation of $\left(\alpha_{1}, \ldots, \alpha_{l}, i\right)(i=1, \ldots, r-1)$. (Notice that each term of (11) can be grouped into pairs $w_{\alpha \beta} \bar{w}_{\gamma \delta}$ in two ways: (i) each pair belongs to the same row, (ii) each pair belongs to the same column.) Since $z_{1 r}=w_{1 r}$, (11) holds for $p=1$. Now assume (11) for $p-1,1 \leqslant r \leqslant m$, and prove for $p$. Upon expanding

$$
\bar{U}_{p}\binom{1, \ldots,[i], r}{1, \ldots, i}
$$

(given by (7)) and multiplying out the resulting factors, we find that its general term is

$$
z_{\lambda_{1} \alpha_{1}} \bar{z}_{\lambda_{1} \beta_{1}} \ldots z_{\lambda_{s} \alpha_{s}} \bar{z}_{\lambda_{s} \beta_{\alpha_{s}}} z_{\lambda_{s+1} r}{\overline{\overline{ }} \lambda_{s+1} \beta_{i}}
$$

where $\lambda_{\alpha}(\alpha=1, \ldots, s+1)$ takes on values from $1,2, \ldots, p-1 ; \alpha_{1}, \ldots, \alpha_{s}$ is a subset of $1, \ldots, i-1$ and

$$
\beta_{\alpha_{1}}, \ldots, \beta_{\alpha_{s}}, \beta_{i}{ }^{\prime}
$$

is a permutation of $\alpha_{1}, \ldots, \alpha_{s}, i$. Thus the general term of $z_{p r}$ would be $w_{p r}$ times

$$
\begin{gathered}
c_{p i}^{\prime} z_{\lambda_{1} \alpha_{1}}{\overline{\lambda_{1} \beta \alpha_{1}}} \ldots z_{\lambda_{s} \alpha_{s}}{\overline{z_{\lambda} s \alpha_{s}}}^{z_{\lambda_{s}+1}} \overline{\bar{z}}_{\lambda_{s}+1 \beta_{i}} w_{p i} \bar{w}_{p r}, \\
c_{p i}{ }^{\prime}=c_{p i} / w_{p r} \bar{w}_{p r} . \text { Replace } \\
z_{\lambda_{1} \alpha_{1} \lambda} \bar{z}_{\lambda_{1} \beta \alpha_{1}} \ldots \text { by } w_{\lambda_{1} \alpha_{1}} \bar{w}_{\lambda_{1} \beta \alpha_{1}} \ldots
\end{gathered}
$$

times a factor which by induction already has the required form and (11) follows.

Now consider $z_{p r}{ }^{s_{p r}}$. From (11) we see that except for the first factor $w_{p r}{ }^{s_{p r}}$ if $w_{\nu j}{ }^{\sigma}$ occurs, then a factor $\bar{w}_{\nu k}{ }^{\sigma}$ also occurs and if $w_{i \mu}{ }^{\tau}$ appears, then $\bar{w}_{l \mu}{ }^{\tau}$ also appears. Consequently in the expression for $z_{p r}{ }^{s_{p r}}$ for each $\nu(\nu=1, \ldots, p-1)$ the sum of the exponents of the factors $w_{\nu j}(j=1, \ldots, m)$ equals the sum of the exponents of the $\bar{w}_{\nu k}(k=1, \ldots, m)$ and the sum of the exponents of $w_{p j}$ equals the sum of the exponents of $\bar{w}_{p k}$ increased by $s_{p r}$. Similarly for the columns. Thus $P$ can be expressed in the form

$$
\begin{equation*}
P(z)=P_{0}(w \bar{w}) \prod w_{j k}^{s_{j k}}=P_{0}(w \bar{w}) P(w), \tag{12}
\end{equation*}
$$

where $P_{0}(w \bar{w})$ contributes the same exponents to the sum of the elements in the $\nu$ th row of $w$ and of $\bar{w}$ and similarly for the columns of $w$ and $\bar{w}$. An analogous expression holds for $Q$.

In $(P, Q)$ replace $P$ and $Q$ by (12). Since $\left\{w_{\nu_{k}}{ }^{s_{\nu k}}\right\}(k=1, \ldots, m)$ forms an orthogonal set on $w_{\nu} w^{*}{ }_{\nu}<1$ (2), $(P, Q) \neq 0$ if and only if for each $k$ the exponent of $w_{\nu k}$ equals the exponent of $\bar{w}_{\nu k}$. Consequently if we sum the exponents of the $\nu$ th row, owing to the form of $P_{0}(w \bar{w}), Q_{0}(w \bar{w})$ we obtain the first of equations (1.8) and summing the exponents of the $\mu$ th column the second of equations (1.8). Thus the theorem is proved.
5. CONS. Orthogonal development. A CONS is constructed from the set of powers $\{P(z)\}$ as follows. Let $\alpha=\left(\alpha_{1}, \ldots, \alpha_{m+n}\right)$ be a set of non-negative integers with $\sum_{j=1}^{n} \alpha_{j}=\sum_{k=1}^{m} \alpha_{k+n}=p$. The powers of the set $S(\alpha)=$ $S\left(\alpha_{1}, \ldots, \alpha_{m+n}\right)$ such that

$$
\sum_{k} s_{\nu k}=\alpha_{\nu}, \quad \sum_{j} s_{j \mu}=\alpha_{\mu+n}
$$

need not be orthogonal to each other. (There exist sets $S(\alpha)$ whose members are not all orthogonal to each other-for example, in the 2 by 2 case the elements $z_{11} z_{22}, z_{12} z_{21}$ are not orthogonal (12).) However if $P \in S(\alpha), Q \in S(\beta)$, where $\alpha_{j} \neq \beta_{j}$ for some $j$, then by Theorem $3.1(P, Q)=0$. We order the elements of the set $S(\alpha)$ in some convenient manner into a sequence $P_{0}, P_{1}, \ldots$, $P_{p(\alpha)}$. An ONS is constructed from these elements by the Gram-Schmidt formulas

$$
\phi_{\nu}^{(p)}(z)=\operatorname{det}\left(\begin{array}{c}
P_{0} \ldots P_{\nu}  \tag{13}\\
a_{00} \ldots a_{\nu 0} \\
\ldots \\
a_{0, \nu-1} \ldots a_{\nu, \nu-1}
\end{array}\right) /\left(D_{\nu-1} D_{\nu}\right)^{\frac{1}{2}},
$$

where

$$
\begin{aligned}
& D_{\mu}=\operatorname{det}\left(a_{\alpha \beta}\right) \quad(0 \leqslant \alpha, \beta \leqslant \mu ; \mu=\nu-1 \text { or } \nu, \nu \neq 0), D_{-1}=1 . \\
& a_{i j}=\left(P_{i}, P_{j}\right),(0 \leqslant i, j \leqslant \nu) .
\end{aligned}
$$

Now order the system $\left\{\phi_{\nu}{ }^{(p)}\right\}$ into a sequence $\phi_{1}, \phi_{2}, \ldots$ The orthogonal development of any $f \in L^{2}$ with respect to the ONS is

$$
\begin{equation*}
\sum a_{q} \phi_{q}, \tag{14}
\end{equation*}
$$

where $a_{q}$ are the Fourier coefficients, $\left(f, \phi_{q}\right)$, of $f$. From Bergman's theory (2) it is known that (14) converges absolutely and continuously to $f$ on $D$ (continuous convergence means that the series converges uniformly on any compact set contained in $D$ ).

## §4. Applications of the CONS.

1. Abelian and Tauberian theorems. Let $\left\{a_{q}\right\}$ be an arbitrary sequence of numbers and consider the behaviour of (3.14) as $z \in D$ approaches a point $u_{0} \in B$. In particular let $u_{0}=[I, 0], I=I^{(n)}$, and $z \rightarrow u_{0}$ along the set of points $[r, 0]$ where $r$ is the diagonal matrix $\left[r_{1}, \ldots, r_{n}\right], 0 \leqslant r_{j}<1$. When $z=[r, 0]$ it is seen that $P_{\nu}$ of the set $S(\alpha)$ is either equal to 0 or to

$$
\Pi r r^{p}
$$

in case $\alpha_{j}=\alpha_{j+n}=p_{j}$ for $j=1, \ldots, n$ and $\alpha_{j+n}=0$ for $j>n$. Consequently for a fixed set $\alpha$ either all $P_{\nu}$ are zero or there is one $P_{\nu}$ different from zero in case $\alpha_{j}=\alpha_{j+n}$ for $j=1, \ldots, n$ and $\alpha_{j+n}=0$ for $j>n$. Order the elements of the set $S(\alpha)$ so that the non-zero term is the last term of the set. Then in (3.13) only $\phi_{t}{ }^{(p)}(z), t=p(\alpha)$, is different from zero when $z=[r, 0]$ and

$$
\phi_{t}^{(p)}(z)=\left[D_{t-1} / D_{t}\right]^{\frac{1}{2}} P_{t}(z)=\left[D_{t-1} / D_{t}\right]^{\frac{1}{2}} r_{1}^{p_{1}} \ldots r_{n}^{p_{n}}
$$

Thus (3.14) reduces to a multiple power series. Let this series be summed by the usual method for power series. Then

$$
\begin{align*}
S(r) & =\sum_{p_{1}=0}^{\infty} \ldots \sum_{p_{n}=0}^{\infty} c_{p_{1} \ldots p_{n}} r_{1}^{p_{1}} \ldots r_{n}^{p_{n}},  \tag{1}\\
c_{p_{1} \ldots p_{n}} & =a_{q} /\left(D_{t-1} / D_{t}\right)^{\frac{1}{2}},
\end{align*}
$$

where $\phi_{t}{ }^{(p)}=\phi_{q}$ is the ordering of the ONS $\left\{\phi_{t}{ }^{(p)}\right\}$ into a simple sequence. Let

$$
S_{q_{1} \ldots q_{n}}
$$

be the partial sum of $S(r)$ :

$$
S_{q_{1} \ldots q_{n}}=\sum_{p_{1}=0}^{q_{1}} \ldots \sum_{p_{n}=0}^{q_{n}} c_{p_{1} \ldots p_{n}} r_{1}^{p_{1}} \ldots r_{n}^{p_{n}} .
$$

The following Abelian theorem is valid:
Theorem 4.1. If $S(I)$ exists and

$$
\left|S_{q_{1}} \cdots q_{n}\right|<C,
$$

where $C$ is an absolute constant, then $S(r)$ is uniformly convergent for $0 \leqslant r_{j}<1$ and $\lim _{r \rightarrow I} S(r)=S(I)$. ( $r \rightarrow I$ means $[r, 0] \rightarrow[I, 0]$.) See (4) for a proof.

Also a Tauberian theorem proved by Knopp (7) for double series may be extended to multiple series.

Theorem 4.2 (Tauberian theorem). Let the series $S(r)$ converge for each $[r, 0] \in D, r=\left[r_{1}, \ldots, r_{n}\right]$, and for these $r$ let $|S(r)| \leqslant K$, where $K$ is an absolute constant. If

$$
\begin{equation*}
\left|c_{p_{1} \ldots p_{n}}\right|\left(p_{1}^{2}+\ldots+p_{n}^{2}\right)^{\frac{1}{2} n}<M<\infty, \tag{2}
\end{equation*}
$$

then $\lim _{r \rightarrow I} S(r)=S$ implies $S(I)=S$, that is,

$$
\begin{equation*}
\sum_{p_{1}=0}^{\infty} \ldots \sum_{p_{n}=0}^{\infty} c_{p_{1} \ldots p_{n}}=S \tag{3}
\end{equation*}
$$

In order to prove Theorem 4.2, Theorems 3 of $\S 3$ and the proofs in $\S 4$ of Knopp's paper must be proved for $n$-fold series ( $n \geqslant 3$ ). Using condition (2) Theorem 3 has been extended to multiple series in (18). Also by means of (2) the proofs in $\S 4$ follow for multiple series. In addition see (13) where a similar condition is utilized for multiple series summed spherically.

On the other hand if we let $z \rightarrow[I, 0]$ along the set $[\rho I, 0], 0 \leqslant \rho<1$, and sum series (3.14) by diagonals:

$$
\begin{equation*}
S_{0}(\rho)=\sum_{p=0}^{\infty} b_{p} \rho^{p} \tag{4}
\end{equation*}
$$

where

$$
\begin{equation*}
b_{p}=\sum_{p_{1}+\ldots+p_{n}=p} c_{p_{1} \ldots p_{n}}, \tag{4a}
\end{equation*}
$$

then (4) is a simple series and the boundedness conditions on $S_{q_{1}} \ldots q_{n}$ and $S(r)$ can be omitted in Theorems 4.1 and 4.2 (17). Abel's theorem reads if $S_{0}(I)$ exists, then $S_{0}(\rho I)$ is uniformly convergent for $0 \leqslant \rho<1$ and $\lim _{\rho \rightarrow 1} S_{0}(\rho I)$ $=S_{0}(I)$ and Theorem 4.2 becomes

Let the series $S_{0}(\rho I)$ converge for each $[\rho I, 0] \in D,(0 \leqslant \rho<1)$. If $b_{p}=0(1 / p)$, then $\lim _{\rho \rightarrow 1} S_{0}(\rho I)=S_{0}$ implies $S_{0}(I)=S_{0}$.
2. Cauchy's inequality and mean value theorem. In the next paragraphs it is convenient to group the elements of the CONS $\left\{\phi_{\nu}\right\}$ of same degree, hence, let

$$
\phi_{1}^{(p)}, \ldots, \phi_{M_{p}}^{(p)}
$$

be the terms of degree $p$. Then for any $f \in L^{2}$ on $D$

$$
\begin{equation*}
f(z)=\sum_{p=0}^{\infty} \sum_{j=1}^{M_{p}} a_{j}^{(p)} \phi_{j}^{(p)}(z), \tag{5}
\end{equation*}
$$

where the convergence is continuous and absolute for $z \in D$. Multiply (5) by $\bar{f}$ and integrate over the domain

$$
\begin{equation*}
D_{\rho}=\left[z \mid \rho^{2} I-z z^{*}>0,0<\rho<1\right] . \tag{6}
\end{equation*}
$$

Clearly $D_{\rho} \nsubseteq D$ for $0<\rho<1$. Since the convergence of (5) is uniform and absolute on $\bar{D}_{\rho}$
(7) $I(\rho)=V_{\rho}^{-1} \int_{D p}|f|^{2} d W_{z}=V_{\rho}^{-1} \sum_{p=0}^{\infty} \sum_{j=1}^{M_{p}} \sum_{q=0}^{\infty} \sum_{k=1}^{M q} a_{j}^{(p)} \bar{a}_{k}^{(q)} \int_{D_{p}} \phi_{j}^{(p)} \bar{\phi}_{k}^{(q)} d W_{z}$
( $V_{\rho}$ being the volume of $D$ ). In the integral

$$
I=\int_{D p} \phi_{j}^{(p)} \bar{\phi}_{k}^{(q)} d W_{z}
$$

set $z=\rho w$. Then $D_{\rho} \rightarrow D, d W_{z} \rightarrow \rho^{2 m n} d W_{w}$ (cf. §3.1) and

$$
\phi_{j}^{(p)}(z)=\sum \Sigma \alpha_{j k=p} \text { constant } \prod z_{j k}^{\alpha_{j k}} \rightarrow \rho^{p} \phi_{j}^{(p)}(w) .
$$

Hence

$$
I=\rho^{2 m n+p+q}\left(\phi_{j}^{(p)}, \phi_{k}^{(q)}\right)=\rho^{2 m n+p+q} \delta_{p q} \delta_{j k} .
$$

Also $V_{\rho}=\rho^{2 m n} V_{0}$ where $V_{0}$ is the volume of $D$. Thus (7) becomes

$$
\begin{equation*}
I(\rho)=\frac{1}{V_{\rho}} \int_{D p}|f|^{2} d W=\frac{1}{V_{0}} \sum_{p=0}^{\infty} \rho^{2 p} \sum_{j=0}^{M_{p}}\left|a_{j}^{(p)}\right|^{2} . \tag{8}
\end{equation*}
$$

From (8)

$$
\frac{1}{V_{0}} \sum_{j=0}^{M_{p}}\left|a_{j}^{(p)}\right|^{2} \leqslant\left(1 / \rho^{2 p}\right) \max _{z \in D p}|f(z)|^{2}
$$

Now according to a theorem proved by Hua (6) if $f$ is analytic on the closed circular domain $\bar{D}_{\rho}$, then $f$ attains its maximum modulus on the circular manifold

$$
\begin{equation*}
B_{\rho}=\left[z \mid z z^{*}=\rho^{2} I\right] . \tag{9}
\end{equation*}
$$

Hence we get the Cauchy inequality:

$$
\begin{equation*}
\frac{1}{V} \sum_{j=0}^{M_{p}}\left|a_{j}^{(p)}\right|^{2} \leqslant\left(1 / \rho^{2 p}\right) \max _{z \in B p}|f(z)|^{2} \tag{10}
\end{equation*}
$$

Theorem 4.3 (mean value theorem). $I(\rho)$ defined by (7) is a monotone increasing function of $\rho(0<\rho<1)$ and $\log I(\rho)$ is a convex function of $\log \rho$.

Proof. The monotonicity of $I(\rho)$ is obvious from (8). The proof of convexity is the same as the proof in ( $\mathbf{1 7}, \mathrm{p} .174$ ) for the one variable case.
3. A mean value theorem over $B_{\rho}$.

Theorem 4.4. In the case $n=m$ the integral

$$
\begin{equation*}
I_{1}(\rho)=\int_{B p}|f|^{2} d V \tag{11}
\end{equation*}
$$

is a monotone increasing function of $\rho$ and $\log I_{1}(\rho)$ is a convex function of $\log \rho$.
Proof. Hua (6) has shown how to construct a set $\left\{\psi_{\nu}\right\}$ orthonormalized with respect to the inner product

$$
\left(\psi_{\nu}, \psi_{\mu}\right)_{B}=\int_{B} \psi_{\nu} \bar{\psi}_{\mu} d V
$$

from the CONS $\left\{\phi_{\nu}\right\}$ as follows. Since $B$ is a circular space $\left(\phi_{j}{ }^{(p)}, \phi_{k}{ }^{(q)}\right)_{B}=0$ if $p \neq q$. Define the vector

$$
z_{p}=\left(\phi_{1}^{(p)}, \ldots, \phi_{M p_{p}}^{(p)}\right)
$$

Then

$$
\left(z_{p}^{\prime}, z_{p}\right)_{B}=\left(\left(\phi_{j}^{(p)}, \phi_{k}^{(p)}\right)_{B}\right)=K_{p}
$$

is a matrix of constants. Since $z_{p}{ }^{\prime} \bar{z}_{p}>0$ if

$$
z_{p} z_{p}^{*}=\left|\phi_{1}^{(p)}\right|^{2}+\ldots+\left|\phi_{M_{p}}^{(p)}\right|^{2}
$$

is positive, $K_{p}>0$ and there exists a unitary matrix $U$ such that $U^{*} K_{p} U=\Lambda$, where $\Lambda$ is a diagonal matrix with positive elements on the diagonal. Now $\left\{y_{p}\right\}$, defined by

$$
y_{p}=z_{p} \bar{U}=\left(\theta_{1}^{(p)}, \ldots, \theta_{M_{p}}^{(p)}\right),
$$

is a CONS on $D$ if $\left\{z_{p}\right\}$ is, since

$$
\left(\left(z_{p} \bar{U}\right)^{\prime}, z_{p} U\right)=U^{*}\left(z_{p}^{\prime}, z_{p}\right) U=U^{*}\left(\left(\phi_{j}^{(p)}, \phi_{k}^{(p)}\right)\right) U=U^{*} I U=I
$$

Let

$$
\psi_{j}^{(p)}=\theta_{j}^{(p)} /\left\|\theta_{j}^{(p)}\right\|_{B}
$$

Then $\left\{\psi_{j}{ }^{(p)}\right\}$ is an ONS with respect to integration over $B$ and the orthogonal development (5) of $f \in L^{2}$ can be written as

$$
\begin{equation*}
f(z)=\sum_{p=0}^{\infty} \sum_{j=1}^{M_{p}} b_{j}^{(p)} \psi_{j}^{(p)} \tag{12}
\end{equation*}
$$

where

$$
b_{j}^{(p)}=\left(f, \psi_{j}^{(p)}\right) /\left\|\psi_{j}^{(p)}\right\|^{2} .
$$

Multiply (12) by $\bar{f}$ and integrate over $B_{\rho}$. By a procedure similar to that in paragraph 3 we obtain the formula

$$
\begin{equation*}
I_{1}(\rho)=\int_{B_{\rho}}|f|^{2} d V=\sum_{p=0}^{\infty} \rho^{2 p} \sum_{j=1}^{M_{p}}\left|b_{j}^{(p)}\right|^{2} \tag{13}
\end{equation*}
$$

from which the conclusions of the theorem follow.
Note added in Proof. Recently it has come to my attention that Hua and Look (21) have proved that $F(z) \rightarrow f\left(u_{o}\right)$ as $z \rightarrow u_{o}$ in any manner. Further for continuous $f$ on $B$, the solution $F$ given by (3) is unique. They also consider Abel summability for continuous functions of the unitary group.

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