# PLURIHARMONIC SYMBOLS OF COMMUTING TOEPLITZ OPERATORS 

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## 1. Introduction and Results

Our setting throughout the paper is the unit ball $B_{n}$ of the complex $n$-space $\mathrm{C}^{n}$; dimension $n$ is fixed and thus we usually write $B=B_{n}$ unless otherwise specified. The Bergman space $A^{2}(B)$ is the closed subspace of $L^{2}(B)=L^{2}(B, V)$ consisting of holomorphic functions where $V$ denotes the volume measure on $B$ normalized to have total mass 1 . For $u \in L^{\infty}(B)$, the Toeplitz operator $T_{u}$ with symbol $u$ is the bounded linear operator on $A^{2}(B)$ defined by $T_{u}(f)=P(u f)$ where $P$ denotes the orthogonal projection of $L^{2}(B)$ onto $A^{2}(B)$. The projection $P$ is the well-known Bergman projection which can be explicitly written as follows:

$$
P(\psi)(z)=\int_{B} \frac{\psi(w)}{(1-\langle z, w\rangle)^{n+1}} d V(w) \quad(z \in B)
$$

for functions $\psi \in L^{2}(B)$. Here $\langle$,$\rangle is the ordinary Hermitian inner product$ on $\mathbf{C}^{n}$. See [7, Chapters 3 and 7] for more information on the projection $P$.

In one dimensional case, Axler and C̆učković [3] has recently obtained a complete description of harmonic symbols of commuting Toeplitz operators: if two Toeplitz operators with harmonic symbols commute, then either both symbols are holomorphic, or both symbols are antiholomorphic, or a nontrivial linear combination of symbols is constant (the converse is also true and trivial). Trying to generalize this characterization to the ball, one may naturally think of pluriharmonic symbols. A function $u \in C^{2}(B)$ is said to be pluriharmonic if its restriction to an arbitrary complex line that intersects the ball is harmonic as a function of single complex variable. As is well known, a real-valued function on $B$ is pluriharmonic if and only if it is the real part of a holomorphic function on $B$. Hence every pluriharmonic function on $B$ can be expressed, uniquely up to an additive constant, as the sum of a holomorphic function and an antiholomorphic function.

[^0]In the present paper we consider the same problem of characterizing pluriharmonic symbols of commuting Toeplitz operators on the ball. Our first result is a necessary condition in terms of $\mathscr{M}$-harmonicity (see Section 2 for relevant definitions) for such symbols.

Theorem 1. Let $f, g, h$, and $k$ be holomorphic functions on $B$ such that $f+\bar{g}$ and $h+\bar{k}$ are pluriharmonic symbols of two commuting Toeplitz operators on $A^{2}(B)$. Then $\overline{f k}-h \bar{g}$ is $\mathscr{M}$-harmonic on $B$.

The proof in [3] shows that the converse of Theorem 1 is also true in one dimensional case. Unfortunately, we were not able to prove or disprove the converse of Theorem 1 on the ball in general. However, Theorem 1 is enough to produce a simple characterization in case one of symbols is holomorphic (or antiholomorphic which amounts to considering adjoint operators). Its proof will make use of a recent characterization (see Proposition 7) of Ahern and Rudin [2] on $\mathscr{M}$-harmonic products.

Theorem 2. Suppose that $u$ and $v$ are pluriharmonic symbols of two commuting Toeplitz operators on $A^{2}(B)$. If $u$ is nonconstant and holomorphic, then $v$ must be holomorphic.

Recall that a bounded linear operator on a Hilbert space is called normal if it commutes with its adjoint operator. Since the adjoint operator of the Toeplitz operator with symbol $u$ is the Toeplitz operator with symbol $\bar{u}$, the following is an immediate consequence of Theorem 2 whose proof is therefore omitted.

Corollary 3. The Toeplitz operator with holomorphic symbol $u$ is normal on $A^{2}(B)$ if and only if $u$ is constant.

In Section 2 we collect some facts about $\mathscr{M}$-harmonic functions which are needed in Section 3 where we prove Theorems 1 and 2. In Section 4 we conclude the paper with some remarks and discussions related to the converse of Theorem 1 and a possible pluriharmonic version of Corollary 3.

## 2. $\mathscr{M}$-Harmonic Functions

For $z, w \in B, z \neq 0$, define

$$
\varphi_{z}(w)=\frac{z-|z|^{-2}\langle w, z\rangle z-\sqrt{1-|z|^{2}}\left(w-|z|^{-2}\langle w, z\rangle z\right)}{1-\langle w, z\rangle}
$$

and $\varphi_{0}(w)=-w$. Then $\varphi_{z} \in \mathscr{M}$, the group of all automorphisms
(= biholomorphic self-maps) of $B$. Furthermore, each $\varphi \in \mathscr{M}$ has a unique representation $\varphi=U \circ \varphi_{z}$ for some $z \in B$ and unitary operator $U$ on $\mathbf{C}^{n}$. For $u \in C^{2}(B)$ and $z \in B$, we define

$$
(\tilde{\Delta} u)(z)=\Delta\left(u \circ \varphi_{z}\right)(0)
$$

where $\Delta$ denotes the ordinary Laplacian. The operator $\tilde{\Delta}$ is called the invariant Laplacian because it commutes with automorphisms of $B$ in the sense that $\tilde{\Delta}(u \circ \varphi)=(\tilde{\Delta} u) \circ \varphi$ for $\varphi \in \mathscr{M}$. We say that a function $u \in C^{2}(B)$ is $\mathscr{M}$-harmonic on $B$ if it is annihilated on $B$ by $\tilde{\Delta}$. One can easily see that $\mathscr{M}$-harmonic functions are precisely harmonic ones in one dimensional case. As is the case for harmonic functions, $\mathscr{M}$-harmonic functions are characterized by a certain mean value property (see [7, Chapter 4]): a function $u \in C(B)$ is $\mathscr{M}$-harmonic on $B$ if and only if

$$
(u \circ \varphi)(0)=\int_{S}(u \circ \varphi)(r \zeta) d \sigma(\zeta) \quad(0 \leq r<1)
$$

for every $\varphi \in \mathscr{M}$. Here $\sigma$ denotes the rotation invariant probability measure on the unit sphere $S$, the boundary of $B$. This is the so-called invariant mean value property. The following area version of this invariant mean value property also gives a characterization of $\mathscr{M}$-harmonicity of functions continuous up to the boundary (see [7, Proposition 13.4.4]): a function $u \in C(\bar{B})$ is $\mathscr{M}$-harmonic on $B$ if and only if

$$
(u \circ \varphi)(0)=\int_{B}(u \circ \varphi) d V
$$

for every $\varphi \in \mathscr{M}$.
The key step to our proof of Theorem 1 is adapted from that of [3]. That is, we will use a slight variant of the characterization of $\mathscr{M}$-harmonicity given by the area version of invariant mean value property. To state it, let us introduce some notations. We associate with each $v \in C(B)$ its so-called radialization $\mathscr{A}(v)$ defined by the formula

$$
\mathscr{A}(v)(z)=\int_{\mathscr{U}}(v \circ U)(z) d U \quad(z \in B)
$$

where $d U$ denotes the Haar measure on the group $\mathscr{U}$ of all unitary operators on $\mathbf{C}^{n}$. Using Proposition 1.4.7 of [7], one can easily verify that

$$
\mathscr{A}(v)(z)=\int_{S} v(|z| \zeta) d \sigma(\zeta) \quad(z \in B)
$$

and hence $\mathscr{A}(v)$ is indeed a radial function on $B$. We write $\mathscr{A}(v) \in C(\bar{B})$ if $\mathscr{A}(v)$ has a continuous extension up to the boundary.

## Proposition 4. Suppose that $u \in C(B) \cap L^{1}(B)$. Then $u$ is $\mathscr{M}$-harmonic on B if and only if

$$
\begin{equation*}
\int_{B}(u \circ \varphi) d V=(u \circ \varphi)(0) \tag{1}
\end{equation*}
$$

and

$$
\begin{equation*}
\mathscr{A}(u \circ \varphi) \in C(\bar{B}) \tag{2}
\end{equation*}
$$

for every $\varphi \in \mathscr{M}$.
Proof. We first prove the easy direction. Suppose that $u$ is $\mathscr{M}$-harmonic on $B$. Let $\varphi \in \mathscr{M}$. By the invariant mean value property, we have

$$
(u \circ \varphi)(0)=\int_{S}(u \circ \varphi)(r \zeta) d \sigma(\zeta)
$$

for every $r \in[0,1)$. Integrating in polar coordinates, we have (1). The above also shows that $\mathscr{A}(u \circ \varphi)$ is constant on $B$, with value $(u \circ \varphi)(0)$, and therefore (2) holds.

To prove the other direction (which we need for the proof of Theorem 1), suppose that (1) and (2) hold. Let $\varphi \in \mathscr{M}$ and put $v=\mathscr{A}(u \circ \varphi)$. We first show that $v$ is $\mathscr{M}$-harmonic on $B$. Since $v \in C(\bar{B})$ by (2), it is sufficient to show the area version of invariant mean value property of $v$. To do this, fix $\psi \in \mathscr{M}$. Then

$$
\begin{equation*}
\int_{B}(v \circ \psi) d V=\int_{B} \int_{\mathscr{U}}\left(u \circ F_{U}\right)(z) d U d V(z) \tag{3}
\end{equation*}
$$

where $F_{U}=\varphi \circ U \circ \psi \in \mathscr{M}$.
For a fixed unitary operator $U$ on $\mathbf{C}^{n}$, consider the inverse mapping $G_{U} \in \mathscr{M}$ of $F_{U}$ and put $a=F_{U}(0)=(\varphi \circ U \circ \psi)(0)$. Then, since $\left|\varphi^{-1}(0)\right|=$ $|\varphi(0)|$, we have ([7, Theorem 2.2.5])

$$
1-|a|^{2}=\frac{\left(1-|\varphi(0)|^{2}\right)\left(1-|\psi(0)|^{2}\right)}{\left|1-\left\langle\varphi^{-1}(0),(U \circ \psi)(0)\right\rangle\right|^{2}} \geq\left(1-|\varphi(0)|^{2}\right)\left(1-|\psi(0)|^{2}\right)
$$

On the other hand, we have [7, Theorem 2.2.6]

$$
J_{R} G_{U}(w)=\left(\frac{1-|a|^{2}}{|1-\langle w, a\rangle|^{2}}\right)^{n+1} \leq\left(\frac{4}{1-|a|^{2}}\right)^{n+1} \quad(w \in B)
$$

where $J_{R} G_{U}(w)$ denotes the real Jacobian of $G_{U}$ at $w \in B$. It follows that the function $J_{R} G_{U}$ is bounded on $B$ uniformly in $U$. Therefore, since $u \in L^{1}(B)$ by assumption, a change of variables shows that

$$
\int_{\mathscr{U}} \int_{B}\left|u \circ F_{U}\right| d V d U=\int_{\mathscr{U}} \int_{B}|u| J_{R} G_{U} d V d U<\infty
$$

Now one can interchange the order of integrations on the right side of (3) to obtain

$$
\begin{aligned}
\int_{B}(v \circ \psi) d V & =\int_{\mathscr{U}} \int_{B}\left(u \circ F_{U}\right) d V d U \\
& =\int_{\mathscr{U}}\left(u \circ F_{U}\right)(0) d U \\
& =\int_{\mathscr{U}}(u \circ \varphi \circ U)(\psi(0)) d U \\
& =\mathscr{A}(u \circ \varphi)(\psi(0)) \\
& =(v \circ \psi)(0)
\end{aligned}
$$

where the second equality holds by (1). Hence $v$ is $\mathscr{M}$-harmonic on $B$. Since $v$ is radial, the invariant mean value property shows that $v$ is constant. Consequently,

$$
(u \circ \varphi)(0)=v(0)=v(z)=\int_{S}(u \circ \varphi)(|z| \zeta) d \sigma(\zeta) \quad(z \in B)
$$

Since $\varphi \in \mathscr{M}$ is arbitrary, the above shows that $u$ has the invariant mean value property and hence that $u$ is $\mathscr{M}$-harmonic on $B$ as desired.

## 3. Proofs

First, we recall some well known facts on the Hardy space $H^{2}(B)$ consisting of holomorphic functions $f$ on $B$ for which

$$
\sup _{0<r<1} \int_{S}|f(r \zeta)|^{2} d \sigma(\zeta)<\infty
$$

Note that $H^{2}(B) \subset A^{2}(B)$ by an integration in polar coordinates. To each $f \in H^{2}(B)$ corresponds its boundary function $f^{*}$ on $S$ defined by $f^{*}(\zeta)=$ $\lim _{r>1} f(r \zeta)$ for $\sigma$-almost every $\zeta \in S$. In addition, we have $f^{*} \in L^{2}(\sigma)$ and

$$
\lim _{r>1} \int_{S}\left|f(r \zeta)-f^{*}(\zeta)\right|^{2} d \sigma(\zeta)=0
$$

See [7, Chapter 5] for details. One can easily verify by using the above that if $f, g \in H^{2}(B)$, then

$$
\lim _{r \not 11} \int_{S} f(r \zeta) \bar{g}(r \zeta) d \sigma(\zeta)=\int_{S} f^{*} \bar{g}^{*} d \sigma
$$

and hence $\mathscr{A}(f \bar{g}) \in C(\bar{B})$.
Next, before turning to the proof of Theorem 1, we prove a couple of lemmas. For $\varphi \in \mathscr{M}$, let $U_{\varphi}$ denote the linear operator on $A^{2}(B)$ defined by $U_{\varphi} f=(f \circ \varphi) J_{\varphi}$ where $J_{\varphi}$ is the complex Jacobian of $\varphi$ and write $U_{\varphi}^{*}$ for its adjoint operator.

Lemma 5. Let $\varphi \in \mathscr{M}$. Then $U_{\varphi} U_{\varphi}^{*}=U_{\varphi}^{*} U_{\varphi}$ is the identity operator on $A^{2}(B)$.

In other words, the conclusion of the lemma is that $U_{\varphi}$ is unitary on $A^{2}(B)$.

Proof. Since $|J \varphi|^{2}$ is the real Jacobian of $\varphi$, a change of variables yields

$$
\int_{B}|(f \circ \varphi)|^{2}|J \varphi|^{2} d V=\int_{B}|f|^{2} d V
$$

for every $f \in A^{2}(B)$, and hence $U_{\varphi}$ is an isometry of $A^{2}(B)$ into $A^{2}(B)$. Clearly $U_{\varphi^{-1}}$ is the inverse operator for $U_{\varphi}$. An invertible linear isometry on a Hilbert space is a unitary operator (see for example [5, Theorem 12.13]). The proof is complete.

Lemma 6. Let $\varphi \in \mathscr{M}$ and let $u \in L^{\infty}(B)$. Then

$$
U_{\varphi} T_{u} U_{\varphi}^{*}=T_{u \circ \varphi}
$$

Recall that $P$ denotes the Bergman projection of $L^{2}(B)$ onto $A^{2}(B)$.
Proof. Define $V_{\varphi}: L^{2}(B) \rightarrow L^{2}(B)$ by $V_{\varphi} f=(f \circ \varphi) J \varphi$. As in the proof of Lemma $5, V_{\varphi}$ is unitary on $L^{2}(B)$. Since $V_{\varphi}=U_{\varphi}$ when restricted to $A^{2}(B)$,
we see that $V_{\varphi}$ takes $A^{2}(B)$ onto $A^{2}(B)$ and hence

$$
\begin{equation*}
P V_{\varphi}=V_{\varphi} P \tag{4}
\end{equation*}
$$

If $f \in A^{2}(B)$, then we see from (4) that

$$
\begin{aligned}
T_{u \circ \varphi} U_{\varphi} f & =T_{u \circ \varphi}((f \circ h) J \varphi)=P((u \circ \varphi)(f \circ \varphi) J \varphi) \\
& =P\left(V_{\varphi}(u f)\right)=V_{\varphi}(P(u f))=U_{\varphi} T_{u} f
\end{aligned}
$$

Thus $T_{u \circ \varphi} U_{\varphi}=U_{\varphi} T_{u}$, and since $U_{\varphi}$ is unitary by Lemma 5, we have $T_{u \circ \varphi}=$ $U_{\varphi} T_{u} U_{\varphi}^{*}$. The proof is complete.

We are now ready to prove Theorem 1.
Proof of Theorem 1. Let $u=f+\bar{g}$ and $v=h+\bar{k}$. Since $u$ and $v$ are bounded on $B$, functions $f, g, h$, and $k$ must be in $H^{2}(B)$ by an application of the Korányi-Vagi theorem (see [7, Theorem 6.3.1]). In particular, functions $f, g, h$, and $k$ are all in $A^{2}(B)$. Let 1 denote the constant function 1 on $B$. Then we have

$$
T_{u} T_{v} 1=T_{u}(P v)=T_{u}(h+\bar{k}(0))=P(f h+\bar{k}(0) f+h \bar{g}+\bar{g}(0) \bar{k}(0)) .
$$

Note that $\int_{B} F d V=F(0)$ for holomorphic functions $F \in L^{1}(B)$. Since the projection $P$ is orthogonal, it follows that

$$
\begin{align*}
\int_{B} T_{u} T_{v} 1 d V & =\int_{B} f h+\bar{k}(0) f+h \bar{g}+\bar{g}(0) \bar{k}(0) d V  \tag{5}\\
& =f(0) h(0)+f(0) \bar{k}(0)+\bar{g}(0) \bar{k}(0)+\int_{B} h \bar{g} d V
\end{align*}
$$

Similarly,

$$
\begin{equation*}
\int_{B} T_{v} T_{u} 1 d V=f(0) h(0)+h(0) \bar{g}(0)+\bar{g}(0) \bar{k}(0)+\int_{B} f \bar{k} d V \tag{6}
\end{equation*}
$$

Since $T_{u} T_{v}=T_{v} T_{u}$ by assumption, letting $\alpha=f \bar{k}-h \bar{g}$, we have by (5) and (6) that

$$
\begin{equation*}
\int_{B} \alpha d V=\alpha(0) \tag{7}
\end{equation*}
$$

We also have (by a remark at the beginning of this section) that

$$
\begin{equation*}
\mathscr{A}(\alpha) \in C(\bar{B}) \tag{8}
\end{equation*}
$$

Let $\varphi \in \mathscr{M}$. Multiplying both sides of the equation $T_{u} T_{v}=T_{v} T_{u}$ by $U_{\varphi}$ on the left and by $U_{\varphi}^{*}$ on the right, we obtain by Lemma 5 that

$$
U_{\varphi} T_{u} U_{\varphi}^{*} U_{\varphi} T_{v} U_{\varphi}^{*}=U_{\varphi} T_{\nu} U_{\varphi}^{*} U_{\varphi} T_{u} U_{\varphi}^{*}
$$

and therefore by Lemma 6

$$
\begin{equation*}
T_{u \circ \varphi} T_{v \circ \varphi}=T_{v \circ \varphi} T_{u \circ \varphi} \tag{9}
\end{equation*}
$$

Equations (7) and (8) were derived under the assumption that $T_{u} T_{v}=T_{v} T_{u}$. Thus (9) says that (7) and (8) remain valid with $\alpha \circ \varphi$ in place of $\alpha$. That is,

$$
\int_{B}(\alpha \circ \varphi) d V=(\alpha \circ \varphi)(0)
$$

and $\mathscr{A}(\alpha \circ \varphi) \in C(\bar{B})$ for any $\varphi \in \mathscr{M}$. It follows from Proposition 4 that $\alpha$ is $\mathscr{M}$-harmonic on $B$. This completes the proof.

Having proved Theorem 1, we now turn to the proof of Theorem 2 which states that if one of symbols of two commuting Toeplitz operators is nonconstant and holomorphic, then the other one must be also holomorphic. In the proof we apply a consequence of the following recent theorem of Ahern and Rudin [2] on $\mathscr{K}$-harmonic products.

Proposition 7. Let $f$ and $g$ be holomorphic functions such that $f \bar{g}$ is $\mathscr{M}$-harmonic on $B$.
(a) If $n \leq 2$, then either $f$ or $g$ is constant.
(b) If $n \geq 3$, and if both $f$ and $g$ are nonconstant, then there exist an integer $2 \leq m \leq n-1$, a unitary operator $U$ on $\mathbf{C}^{n}$, and entire functions $F$ on $\mathbf{C}^{m-1}$, and $G$ on $\mathbf{C}^{n-m}$, such that

$$
f(U z)=F\left(\frac{z_{2}}{1-z_{1}}, \ldots, \frac{z_{m}}{1-z_{1}}\right), \quad g(U z)=G\left(\frac{z_{m+1}}{1-z_{1}}, \ldots, \frac{z_{n}}{1-z_{1}}\right)
$$

Moreover, $f(B)=F\left(\mathbf{C}^{m-1}\right), \quad g(B)=G\left(\mathbf{C}^{n-m}\right)$, and $(f \bar{g})(B)=\mathbf{C}$ or $\mathbf{C} \backslash\{0\}$.

Combining Proposition 7 with Liouville's theorem, we have the following:
Lemma 8. Let $f$ and $g$ be holomorphic functions such that $f \bar{g}$ is $\mathscr{M}$-harmonic on B. If one of them is bounded on $B$, then either $f$ or $g$ is constant.

Proof of Theorem 2. Write $v=h+\bar{k}$ where $h, k$ are holomorphic on $B$. Then, by Theorem $1, u k$ is $\mathscr{M}$-harmonic on $B$. Since $u$ is bounded and nonconstant on $B$ by assumption, we see from Lemma 8 that $k$ must be constant and hence $v$ is holomorphic on $B$. Conversely, since Toeplitz operators with holomorphic symbols are simply multiplication operators, it is straightforward that two Toeplitz operators with holomorphic symbols commute on $A^{2}(B)$.

## 4. Some related remarks

Throughout the section $f, g, h$, and $k$ denote holomorphic functions on $B$, normalized so that $f(0)=g(0)=h(0)=k(0)=0$ for simplicity. In view of Theorem 1 one may ask (under additional boundedness hypothesis as in Lemma 8 if desired) whether there is any further description of such functions for which

$$
\begin{equation*}
\tilde{\Delta}(f \bar{k})=\tilde{\Delta}(h \bar{g}) \tag{10}
\end{equation*}
$$

Both sides of the above are assumed to be not identically zero; otherwise we are back to Proposition 7. In one dimensional case, it is elementary to verify that condition (10) implies $f=\lambda h$ and $g=\bar{\lambda} k$ for some constant $\lambda$. In higher dimensional cases, we do not know whether the same is true in general. This question can be rephrased as follows: does it follow from (10) that $f \bar{k}-h \bar{g}$ is pluriharmonic? The answer is known to be yes if an additional smoothness condition of certain order, depending on dimension $n$, is satisfied up to the boundary: if a function $u \in C^{n}(\bar{B})$ is $\mathscr{M}$-harmonic on $B$, then $u$ is pluriharmonic on $B$. See [1] or [4]. We also remark in passing that there is in fact a more precise version of this fact ([6]): if $u$ is $\mathscr{M}$-harmonic on $B$ and if the $n$th radial derivative $\mathscr{D}^{n} u$ satisfies the $L^{2}$-growth condition

$$
\left\{\int_{S}\left|\left(\mathscr{D}^{n} u\right)(r \zeta)\right|^{2} d \sigma(\zeta)\right\}^{1 / 2}=o\left(\log \frac{1}{1-r}\right) \quad(r \nearrow 1)
$$

then $u$ is pluriharmonic on $B$. Note that $T_{f+\bar{g}} T_{h+\bar{k}}=T_{h+\bar{k}} T_{f+\bar{g}}$ if and only if $T_{f} T_{\bar{k}}-T_{\bar{k}} T_{f}=T_{h} T_{\bar{g}}-T_{\bar{g}} T_{h}$ for functions $f, g, h$, and $k$ bounded on $B$. Thus, for example, we have the following:

If $f, g, h$, and $k$ are of class $C^{n}$ on $\bar{B}$, and if $T_{f} T_{\bar{k}}-T_{\bar{k}} T_{f}=T_{h} T_{\bar{g}}-T_{\bar{g}} T_{h}$ on $A^{2}(B)$, then $f=\lambda h$ and $g=\bar{\lambda} k$ for some constant $\lambda$.

Trying to obtain a pluriharmonic version of Corollary 3, one is led to a special case of (10) which may be of some independent interest. That is, the
question is now whether the condition

$$
\begin{equation*}
\tilde{\Delta}|f|^{2}=\tilde{\Delta}|g|^{2} \tag{11}
\end{equation*}
$$

implies $f=\lambda g$ for some unimodular constant $\lambda$. We could prove only in some special cases that the answer is yes. Those are included in the rest of the paper with hope that they may serve as a motivation for someone to settle the question in the affirmative or negative direction. We first prove a couple of lemmas.

Lemma 9. Let $\Omega$ be a given connected open subset of $\mathbf{C}^{n}$. If $F_{j}$ and $G_{j}$ ( $1 \leq j \leq m$ ) are holomorphic functions such that $\sum_{j=1}^{m} F_{j} \bar{G}_{j}=0$ on $\Omega$, then $\sum_{j=1}^{m} F_{j}(z) \bar{G}_{j}(w)=0$ for all $z, w \in \Omega$.

Proof. Assume, without loss of generality, that an open ball $\beta$ with center at the origin is contained in $\Omega$. Define

$$
H(z, w)=\sum_{j=1}^{m} F_{j}(z) \bar{G}_{j}(\bar{w}) \quad(z, \bar{w} \in \Omega)
$$

It is sufficient to show that $H(z, \bar{w})=0$ for all $z, w \in \beta$ by real analyticity. Let $L$ be the invertible linear operator on $\mathbf{C}^{n} \times \mathbf{C}^{n}$ defined by

$$
L(z, w)=(z+i w, z-i w)
$$

Then, since $H(z, \bar{z})=0$ for all $z \in \beta$ by hypothesis, we have $H \circ L=0$ on $V \cap\left(\mathbf{R}^{n} \times \mathbf{R}^{n}\right)$ where $V=L^{-1}(\beta \times \beta)$. Note that the function $H \circ L$ is holomorphic on $V$. A consideration of Taylor coefficients therefore shows that $H \circ L$ vanishes on $V$. In other words, $H=0$ on $\beta \times \beta$, completing the proof.

Lemma 10. Let $\Omega$ be a given connected open subset of $\mathbf{C}^{n}$. If $F_{j}$ and $G_{j}$ $(1 \leq j \leq m)$ are holomorphic functions such that $\sum_{j=1}^{m}\left|F_{j}\right|^{2}=\sum_{j=1}^{m}\left|G_{j}\right|^{2}$ on $\Omega$, then there is a unitary operator $U$ on $\mathbf{C}^{m}$ such that $\left(F_{1}, \ldots, F_{m}\right)=$ $U \circ\left(G_{1}, \ldots, G_{m}\right)$ on $\Omega$.

Proof. The lemma is trivial if $m=1$. To proceed by induction on $m$, let $m>1$ and suppose that the lemma is proved for $m-1$. Let $F=\left(F_{1}, \ldots, F_{m}\right)$ and $G=\left(G_{1}, \ldots, G_{m}\right)$. We may assume that $\Omega$ contains the origin. We may further assume that $|F(0)|=|G(0)|=1$. Pick unitary operators $U_{1}$ and $U_{2}$ on $\mathbf{C}^{m}$ such that $U_{1}(F(0))=U_{2}(G(0))=(1,0, \ldots, 0)$. Let

$$
U_{1} \circ F=\left(f_{1}, \ldots, f_{m}\right) \text { and } U_{2} \circ G=\left(g_{1}, \ldots, g_{m}\right)
$$

Then we have $\sum_{j=1}^{m} f_{j} \bar{f}_{j}=\sum_{j=1}^{m} g_{j} \bar{g}_{j}$ on $\Omega$ and hence, by Lemma 9,

$$
\sum_{j=1}^{m} f_{j}(z) \bar{f}_{j}(w)=\sum_{j=1}^{m} g_{j}(z) \bar{g}_{j}(w)
$$

for all $z, w \in \Omega$. Taking $w=0$, we obtain $f_{1}=g_{1}$ on $\Omega$. Thus, by induction hypothesis, there exists some unitary operator $U$ on $\mathbf{C}^{m-1}$ such that $\left(f_{2}, \ldots, f_{m}\right)=U \circ\left(g_{2}, \ldots, g_{m}\right)$ on $\Omega$. Now let

$$
U_{3}=\left(\begin{array}{cccc}
1 & 0 & \cdots & 0 \\
0 & & & \\
\vdots & & U & \\
0 & & &
\end{array}\right)
$$

Then $U_{3}$ is a unitary operator on $\mathbf{C}^{m}$ and we have $F=U_{1}^{-1} \circ U_{3} \circ U_{2} \circ G$. The proof is complete.

In what follows, we let $\nabla f=\left(D_{1} f, \ldots, D_{n} f\right)$ and $\mathscr{R} f=\sum_{j=1}^{n} z_{j} D_{j} f$ where $D_{j}$ denotes the differentiation with respect to $z_{j}$-variable. With these notations, equation (11) becomes

$$
\begin{equation*}
|\nabla f|^{2}+|\mathscr{R} g|^{2}=|\nabla g|^{2}+|\mathscr{R} f|^{2} \tag{12}
\end{equation*}
$$

We assert the following:
Suppose that (12) holds on $B_{2}$. If $\nabla f(0)=\nabla g(0) \neq 0$, then $f=g$ on $B_{2}$.
Proof. By (12) and Lemma 10 there is a unitary operator ( $\alpha_{i j}$ ) on $\mathbf{C}^{3}$ such that

$$
\left(\begin{array}{c}
D_{1} f  \tag{13}\\
D_{2} f \\
\mathscr{R} g
\end{array}\right)=\left(\begin{array}{lll}
\alpha_{11} & \alpha_{12} & \alpha_{13} \\
\alpha_{21} & \alpha_{22} & \alpha_{23} \\
\alpha_{31} & \alpha_{32} & \alpha_{33}
\end{array}\right)\left(\begin{array}{c}
D_{1} g \\
D_{2} g \\
\mathscr{R} f
\end{array}\right)
$$

Assume that $\nabla f(0)=\nabla g(0)=(1,0)$ without loss of generality. Then, evaluating both sides of (13) at the origin, one can easily find that $\alpha_{11}=1$ and $\alpha_{12}=\alpha_{13}=\alpha_{21}=\alpha_{31}=0$. It follows that $D_{1} f=D_{1} g$. Hence $D_{2} f-D_{2} g$ does not depend on $z_{1}$-variable. In order to prove $D_{2} f=D_{2} g$ it is therefore sufficient to show that $D_{2} f\left(0, z_{2}\right)=D_{2} g\left(0, z_{2}\right)$ for $\left|z_{2}\right|<1$. Evaluating both sides of (12) at points ( $0, z_{2}$ ), we obtain that $\left|D_{2} f\left(0, z_{2}\right)\right|=\left|D_{2} g\left(0, z_{2}\right)\right|$ and thus there exists a unimodular constant $\lambda$ such that

$$
\begin{equation*}
D_{2} f\left(0, z_{2}\right)=\lambda D_{2} g\left(0, z_{2}\right) \tag{14}
\end{equation*}
$$

for $\left|z_{2}\right|<1$. Assume that both sides of (14) are not identically zero; otherwise we are done. By (13),

$$
z_{2} D_{2} g\left(0, z_{2}\right)=\alpha_{32} D_{2} g\left(0, z_{2}\right)+\alpha_{33} z_{2} D_{2} f\left(0, z_{2}\right)
$$

Insert (14) into the above. A little manipulation yields $\alpha_{32}=\alpha_{23}=0$ and $\alpha_{33}=\bar{\lambda}$. Thus, we have $\mathscr{R} f=\lambda \mathscr{R} g$. Evaluating both sides of this at points ( $z_{1}, 0$ ), we obtain $\lambda=1$. The proof is complete.

We now conclude the paper with another special case:
If $f_{l+1}=0$ and $f_{j} \neq 0$ for some $1 \leq j \leq l$ where $f_{m}$ denotes the mth degree term in the homogeneous expansion of $f$ on $B$, then (12) implies $f=\lambda g$ for some unimodular constant $\lambda$.

Thus, if there were counter examples, then there would be no "gap" in their homogeneous expansions.

Proof. First note that the invariant mean value property of $|f|^{2}-|g|^{2}$ yields

$$
\int_{S}\left|f_{m}\right|^{2} d \sigma=\int_{S}\left|g_{m}\right|^{2} d \sigma \quad(m=1,2, \ldots)
$$

where $g_{m}$ denotes the $m$ th degree term in the homogeneous expansion of $g$ on $B$. Hence $g_{l+1}=0$ and $g_{j} \neq 0$ by hypothesis. Now, by Lemma 10 as before, there exists a unitary operator $U$ on $\mathbf{C}^{n+1}$ such that $(\nabla f, \mathscr{R} g)=$ $U \circ(\nabla g, \mathscr{R} f)$. In particular, there are some vectors $\alpha, \beta \in \mathbf{C}^{n}$ and a constant $\lambda$ with $|\alpha|^{2}+|\lambda|^{2}=|\beta|^{2}+|\lambda|^{2}=1$ such that

$$
\begin{equation*}
\mathscr{R} f=\langle\nabla g, \alpha\rangle+\lambda \mathscr{R} g \quad \text { and } \quad \mathscr{R} g=\langle\nabla f, \beta\rangle+\bar{\lambda} \mathscr{R} f \tag{15}
\end{equation*}
$$

If $|\lambda|=1$, then $\alpha=\beta=0$ and (15) shows $\mathscr{R} f=\lambda \mathscr{R} g$, hence $f=\lambda g$. So, we assume $|\lambda|<1$ and derive a contradiction. Equate terms of the same degree in the homogeneous expansions of both sides of two equations of (15) to obtain

$$
m f_{m}=\left\langle\nabla g_{m+1} m, \alpha\right\rangle+\lambda m g_{m} \quad \text { and } \quad m g_{m}=\left\langle\nabla f_{m+1}, \beta\right\rangle+\bar{\lambda} m f_{m}
$$

so that

$$
m\left(1-|\lambda|^{2}\right) f_{m}=\left\langle\nabla g_{m+1}, \alpha\right\rangle+\lambda\left\langle\nabla f_{m+1}, \beta\right\rangle
$$

for $m=1,2, \ldots$. Since $f_{l+1}=g_{l+1}=0$, the above shows that $f_{m}=g_{m}=0$ for all $1 \leq m \leq l$, which is a contradiction. The proof is complete.

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## References

[1]. P. Ahern and K. Johnson, Differentiability criteria and harmonic functions on $B^{n}$, Proc. Amer. Math. Soc., vol. 89 (1983), pp. 709-712.
[2]. P. Ahern and W. Rudin, $\mathscr{M}$-harmonic products, Indag. Math., vol. 2 (1991), pp. 141-147.
[3]. S. Axler and Z̆. C̆učković, Commuting Toeplitz operators with harmonic symbols, Integral Equations Operator Theory, vol. 14 (1991), pp. 1-11.
[4]. R. Graham, The Dirichlet problem for the Bergman Laplacian, Comm. Partial Differential Equations, vol. 8 (1983), pp. 433-476.
[5]. W. Rudin, Functional analysis, McGraw-Hill, New York, 1973.
[6]. $\qquad$ , A smoothness condition that implies pluriharmonicity, unpublished.
[7]. $\qquad$ , Function theory in the unit ball of $\mathbf{C}^{n}$, Springer-Verlag, New York, 1980.

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