# GEOMETRIC ARVESON-DOUGLAS CONJECTURE FOR THE DRURY-ARVESON SPACE: THE CASE OF ONE-DIMENSIONAL VARIETY 

Jingbo Xia


#### Abstract

We consider a class of analytic subsets $\tilde{M}$ of an open neighborhood of the closed unit ball in $\mathbf{C}^{n}$. Such an $\tilde{M}$ gives rise to a submodule $\mathcal{R}$ and a quotient module $\mathcal{Q}$ of the Drury-Arveson module $H_{n}^{2}$ in $n$ variables. The geometric Arveson-Douglas conjecture predicts that the quotient module $\mathcal{Q}$ is $p$-essentially normal for $p>d=\operatorname{dim}_{\mathbf{C}} \tilde{M}$. We prove this conjecture for the case of dimension $d=1$. In fact, we prove that if $d=1$, then $\mathcal{Q}$ is 1-essentially normal, which is a stronger result than the original prediction.


## 1. Introduction

Let $\mathbf{B}$ denote the open unit ball $\left\{z \in \mathbf{C}^{n}:|z|<1\right\}$ in $\mathbf{C}^{n}$. Throughout the paper, the complex dimension $n$ is always assumed to be greater than or equal to 2. Recall that the Drury-Arveson space $H_{n}^{2}$ is the Hilbert space of analytic functions on $\mathbf{B}$ that has the function

$$
\begin{equation*}
K_{z}(\zeta)=\frac{1}{1-\langle\zeta, z\rangle} \tag{1.1}
\end{equation*}
$$

as its reproducing kernel $[1,16]$. Equivalently, $H_{n}^{2}$ can be described as the Hilbert space of analytic functions on $\mathbf{B}$ where the inner product is given by

$$
\langle h, g\rangle=\sum_{\alpha \in \mathbf{Z}_{+}^{n}} \frac{\alpha!}{|\alpha|!} a_{\alpha} \overline{b_{\alpha}}
$$

for

$$
h(\zeta)=\sum_{\alpha \in \mathbf{Z}_{+}^{n}} a_{\alpha} \zeta^{\alpha} \quad \text { and } \quad g(\zeta)=\sum_{\alpha \in \mathbf{Z}_{+}^{n}} b_{\alpha} \zeta^{\alpha} .
$$

Here and throughout, we follow the standard multi-index notation [26, page 3].
Nowadays, it is common to view $H_{n}^{2}$ as a Hilbert module over the polynomial ring $\mathbf{C}\left[\zeta_{1}, \ldots, \zeta_{n}\right][6,12]$. Thus $H_{n}^{2}$ has submodules and quotient modules.

Suppose that $\mathcal{N}$ is either a submodule or a quotient module of the Drury-Arveson module $H_{n}^{2}$. Let $P_{\mathcal{N}}: H_{n}^{2} \rightarrow \mathcal{N}$ be the orthogonal projection. Then we have the module operators

$$
\mathcal{Z}_{\mathcal{N}, j}=P_{\mathcal{N}} M_{\zeta_{j}} \mid \mathcal{N}, \quad j=1, \ldots, n
$$

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on $\mathcal{N}$. Recall that $\mathcal{N}$ is said to be $p$-essentially normal if all commutators $\left[\mathcal{Z}_{\mathcal{N}, i}, \mathcal{Z}_{\mathcal{N}, j}^{*}\right]$, $1 \leq i, j \leq n$, are in the Schatten class $\mathcal{C}_{p}$.

The famous Arveson Conjecture [2-4] predicts that every graded submodule of $H_{n}^{2} \otimes \mathbf{C}^{r}$ is $p$-essentially normal for $p>n$. To date, the best results on the Arveson Conjecture are due to Guo and K. Wang [20].

In addition to graded submodules, quotient modules of the form $H_{n}^{2} /[I]$, where $I$ is a homogeneous ideal in $\mathbf{C}\left[\zeta_{1}, \ldots, \zeta_{n}\right]$, were also studied in [20]. Motived by the results in [20], Douglas observed in [10] that for quotient modules the essential normality should really be $p>d$, where $d$ is the complex dimension of the variety involved. This more refined conjecture is now called Arveson-Douglas Conjecture. See [5,11,13-15,17,19,20-22,27-29] for the tremendous progress that has been made in this direction.

In this paper we consider a very specific class of submodules and the corresponding quotient modules. Our focus is on the quotient modules, because that is where things become really interesting. Let $\tilde{M}$ be an analytic subset [8] of an open neighborhood of $\overline{\mathbf{B}}$ with $1 \leq \operatorname{dim}_{\mathbf{C}} \tilde{M} \leq n-1$. We will assume that $\tilde{M}$ has no singular points on the sphere $S=\left\{z \in \mathbf{C}^{n}:|z|=1\right\}$ and that $\tilde{M}$ intersects $S$ transversely. Denote $M=\mathbf{B} \cap \tilde{M}$. Then we have a submodule

$$
\mathcal{R}=\left\{f \in H_{n}^{2}: f=0 \text { on } M\right\}
$$

of $H_{n}^{2}$. The corresponding quotient module is

$$
\mathcal{Q}=H_{n}^{2} \ominus \mathcal{R}
$$

In this setting, we have
Geometric Arveson-Douglas Conjecture. The quotient module $\mathcal{Q}$ is $p$-essentially normal for every $p>d=\operatorname{dim}_{\mathbf{C}} \tilde{M}$.

Since the Drury-Arverson module $H_{n}^{2}$ itself is known to be $p$-essentially normal for $p>n$ [1], by a well-known result of Arveson [2], the geometric Arveson-Douglas conjecture implies that the submodule $\mathcal{R}$ is $p$-essentially normal for $p>n$.

The analogous problems for the Bergman module $L_{a}^{2}(\mathbf{B})$ and the Hardy module $H^{2}(S)$ were recently solved $[15,27,28]$. Thus it is logical for us to consider the Arveson-Douglas conjecture for $H_{n}^{2}$. But the case of the Drury-Arverson space $H_{n}^{2}$ poses significant challenges. In fact, it is very easy to describe the main difficulty. Note that the power for the denominator on the right-hand side of (1.1) is only 1 . By contrast, the reproducing kernels for Bergman space $L_{a}^{2}(\mathbf{B})$ and the Hardy space $H^{2}(S)$ are

$$
K_{z}^{\mathrm{Berg}}(\zeta)=\frac{1}{(1-\langle\zeta, z\rangle)^{n+1}} \quad \text { and } \quad K_{z}^{\mathrm{Har}}(\zeta)=\frac{1}{(1-\langle\zeta, z\rangle)^{n}}
$$

which have powers $n+1$ and $n$ for the denominator respectively. Because of the necessary estimates involved, the smaller the power of the reproducing kernel, the harder it is to prove essential normality. Indeed by a comparison of [15,27] with [28], we can already see a significant rise of the level of difficulty when that power is reduced from $n+1$ to $n$.

The challenge we face in this paper is to reduce the power of the reproducing kernel from $n$ all the way to 1 and still prove essential normality. As of this writing, we have managed to overcome this challenge only in the case $d=1$. But the case $d=1$ is also arguably the most interesting one, for in this case we have a stronger result than the prediction of the Arveson-Douglas conjecture:

Theorem 1.1. Suppose that $\operatorname{dim}_{\mathbf{C}} \tilde{M}=1$. Then the quotient module $\mathcal{Q}$ is 1-essentially normal, i.e., every commutator $\left[\mathcal{Z}_{\mathcal{Q}, i}, \mathcal{Z}_{\mathcal{Q}, j}^{*}\right]$ is in the trace class $\mathcal{C}_{1}, i, j \in\{1, \ldots, n\}$.

The rest of the paper is devoted to the long proof of this result. We conclude the Introduction with a brief discussion of the main steps in the proof and the organization of the paper, which should give the reader some idea why the proof is as long as it is.

Our proof of Theorem 1.1 begins with the preliminaries in Section 2. Specifically, in Section 2 we first record the precise definitions of $\tilde{M}, \mathcal{R}, \mathcal{Q}$, etc, and then we collect a number of previously-established results that will be needed in the subsequent sections.

As we have already mentioned, a major difficulty we face is that the power in (1.1) is too small. Our main idea of dealing with this is to increase the power of the denominator by differentiation. Since we only consider the case $d=1$ in the proof of Theorem 1.1, one order of derivative will increase the power of the denominator to 2 , which should be enough based on dimensional considerations. But derivative has to be taken very carefully in the following sense. For $z \in M$, we have $K_{z} \in \mathcal{Q}$, and differentiation modifies the kernel $K_{z}$. We must make sure that the modified kernel still belongs to the quotient module $\mathcal{Q}$. Thus we can only differentiate in the directions tangential to $M$.

To carry out the idea explained above, starting Section 3 we assume $d=1$. We consider a smooth part $\mathcal{M}$ of $\tilde{M}$ near $S$. For each $w \in \mathcal{M}$, let $p_{w}$ be the orthogonal projection of $w$ on the tangent space $T_{w}$. The transversality of $\tilde{M}$ implies that $p_{w} \neq 0$ if $w \in \tilde{M} \cap S$. Thus if $\mathcal{M}$ is a part of $\tilde{M}$ sufficiently near $S$, then $p_{w} \neq 0$ for every $w \in \mathcal{M}$. This gives us a non-vanishing cross section $w \mapsto p_{w}$ of the complex tangent bundle of $\mathcal{M}$. For $w \in \mathcal{M} \cap \mathbf{B}$, the kernel

$$
K_{w, p_{w}}(\zeta)=\frac{\left\langle\zeta, p_{w}\right\rangle}{(1-\langle\zeta, w\rangle)^{2}}
$$

reproduces the derivative in the direction of $p_{w}$. That is, $\left\langle f, K_{w, p_{w}}\right\rangle=\left(\partial_{p_{w}} f\right)(w)$ for $f \in H_{n}^{2}$. Moreover, because $p_{w} \in T_{w}$, we have $K_{w, p_{w}} \in \mathcal{Q}$. Specific to the case $d=1$, we introduce the measure

$$
d \mu(w)=\left(1-|w|^{2}\right) d v_{M}(w)
$$

on $M$. This naturally leads to the operator

$$
T_{1}=\int_{M^{\left(t_{0}\right)}} K_{w, p_{w}} \otimes K_{w, p_{w}} d \mu(w)
$$

on $H_{n}^{2}$, where $M^{\left(t_{0}\right)}$ is a carefully chosen subset of $\mathcal{M} \cap \mathbf{B}$. A major step in the proof of Theorem 1.1 is Theorem 3.5, which says that the spectrum of the positive operator $T_{1}$ does not intersect the interval $(0, c)$ for some $c>0$. Note that even though the cross section
$w \mapsto p_{w}$ is non-vanishing on $\mathcal{M}$, in general it is not analytic. But the condition $d=1$ means that, locally, $w \mapsto p_{w}$ is an analytic cross section multiplied by a scalar function. We will use this fact in the proof of Theorem 3.5.

In addition to $T_{1}$, in Section 4 we introduce the more conventional operator

$$
T_{2}=\int_{M^{\left(t_{0}\right)}} K_{w} \otimes K_{w} d \mu(w)
$$

For our purpose, the operator that really matters is $T=T_{1}+T_{2}$. Theorem 3.5 allows us to show that there is a $c^{\prime}>0$ such that the spectrum of $T$ does not intersect $\left(0, c^{\prime}\right)$, and that $Q$, the orthogonal projection from $H_{n}^{2}$ onto $\mathcal{Q}$, equals the spectral projection of $T$ corresponding to the interval $\left[c^{\prime}, \infty\right)$. In other words, we have a practical control of the orthogonal projection $Q: H_{n}^{2} \rightarrow \mathcal{Q}$ through the operators $T_{1}$ and $T_{2}$, which are given by explicit formulas.

We then introduce a particular Hilbert space $\mathcal{L}$ in Section 4, which can be thought of as a collection of functions on $M^{\left(t_{0}\right)}$ with a particular norm $\|\cdot\|_{\#}$. Let $\mathcal{P}$ be the closure of the analytic polynomials $\mathbf{C}\left[\zeta_{1}, \ldots, \zeta_{n}\right]$ in $\mathcal{L}$. The norm $\|\cdot\|_{\#}$ has the property that

$$
\langle T f, f\rangle=\|f\|_{\#}^{2}
$$

for every $f \in H_{n}^{2}$. This leads to the operator $J$, which take each $f \in H_{n}^{2}$ to the same function $f$ in $\mathcal{P}$. We think of $J$ as restricting each $f \in H_{n}^{2}$ to the set $M^{\left(t_{0}\right)}$. The above identity means that

$$
J^{*} J=T .
$$

Thus it follows from the properties of $T$ that $J: \mathcal{Q} \rightarrow \mathcal{P}$ is invertible. We call $\mathcal{P}$ the range space of the restriction operator $J$. One can think of $\mathcal{P}$ as a representation of the quotient module $\mathcal{Q}$ that is more accessible.

Accordingly, the operators $T_{1}$ and $T_{2}$ also have their representations $\hat{T}_{1}$ and $\hat{T}_{2}$ on $\mathcal{L}$. Individually, the operators $\hat{T}_{1}$ and $\hat{T}_{2}$ are not self-adjoint on $\mathcal{L}$. It is, therefore, something of a miracle that the sum $\hat{T}=\hat{T}_{1}+\hat{T}_{2}$ actually is self-adjoint on $\mathcal{L}$. From this selfadjointness it follows that, with respect to the space decomposition $\mathcal{L}=\mathcal{P} \oplus \mathcal{P}^{\perp}$, we have the operator decomposition $\hat{T}=\tilde{T} \oplus 0$, where $\tilde{T}=J J^{*}$. This means that the orthogonal projection $P: \mathcal{L} \rightarrow \mathcal{P}$ can be expressed in the form $P=\xi(\hat{T})$ for some $\xi \in C_{c}^{\infty}(\mathbf{R})$. This converts the proof of the 1-essential normality in Theorem 1.1 to a problem in terms of commutators and double commutators on $\mathcal{L}$ that are much more accessible than the ones on the Drury-Arveson space $H_{n}^{2}$. But the actual handling of these commutators and double commutators on $\mathcal{L}$ is quite tedious: it takes the work in Sections 5,6 and 7 to finally obtain the 1 -essential normality promised in Theorem 1.1.

There is a major difference between proving essential normality in the case of $H_{n}^{2}$ and the corresponding task in the case of the Bergman space $L_{a}^{2}(\mathbf{B})$ or the Hardy space $H^{2}(S)$. In the two latter cases, the commuting tuple $\left(M_{\zeta_{1}}, \ldots, M_{\zeta_{n}}\right)$ is jointly subnormal, i.e., it extends to a tuple of multiplication operators $\left(M_{\zeta_{1}}, \ldots, M_{\zeta_{n}}\right)$ on an $L^{2}$-space. On the $L^{2}$-space, we have $M_{\zeta_{j}}^{*}=M_{\bar{\zeta}_{j}}, 1 \leq j \leq n$. More to the point, $M_{\zeta_{j}}$ commutes with $M_{\zeta_{i}}^{*}$ on
the $L^{2}$-space, which is a fact heavily involved in [27,28]. In contrast, on the Drury-Arveson space $H_{n}^{2}$, the tuple $\left(M_{\zeta_{1}}, \ldots, M_{\zeta_{n}}\right)$ is known not to be jointly subnormal [1]. This creates an additional obstacle to the proof of essential normality on $H_{n}^{2}$.

Because ( $M_{\zeta_{1}}, \ldots, M_{\zeta_{n}}$ ) is not jointly subnormal on $H_{n}^{2}$, the only way to obtain the desired essential normality for $\mathcal{Q}$ is through the pair of spaces $\mathcal{P} \subset \mathcal{L}$. We have the tuple of multiplication operators $\left(M_{\zeta_{1}}, \ldots, M_{\zeta_{n}}\right)$ on $\mathcal{P}$ with the relation

$$
M_{\zeta_{i}} J=J \mathcal{Z}_{\mathcal{Q}, i}
$$

$1 \leq i \leq n$. The tuple $\left(M_{\zeta_{1}}, \ldots, M_{\zeta_{n}}\right)$ on $\mathcal{P}$ naturally extends to the commuting tuple $\left(\hat{M}_{\zeta_{1}}, \ldots, \hat{M}_{\zeta_{n}}\right)$ on $\mathcal{L}$. On $\mathcal{L}$, we still have $\hat{M}_{\zeta_{i}}^{*} \neq \hat{M}_{\bar{\zeta}_{i}}$, but the difference $\hat{M}_{\zeta_{i}}^{*}-\hat{M}_{\bar{\zeta}_{i}}$ can be computed explicitly. In fact, the handling of this difference is a significant part of Sections 5,6 and 7 . But what matters is the fact that in the end, this approach does lead to a proof of Theorem 1.1.

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## 2. Preliminaries

Although the actual work of this paper only concerns the case $\operatorname{dim}_{\mathbf{C}} \tilde{M}=1$, due to the need to cite a number of existing results, it is necessary to introduce the general technical framework. Much of the material in this section is cited from [28, Section 2].

We begin with the precise definitions of the analytic sets, submodules and quotient modules that we consider in this paper.

Definition 2.1. [8] Let $\Omega$ be a complex manifold. A set $A \subset \Omega$ is called a complex analytic subset of $\Omega$ if for each point $a \in \Omega$ there are a neighborhood $U$ of $a$ and functions $f_{1}, \cdots, f_{N}$ analytic in this neighborhood such that

$$
A \cap U=\left\{z \in U: f_{1}(z)=\cdots=f_{N}(z)=0\right\} .
$$

A point $a \in A$ is called regular if there is a neighborhood $U$ of $a$ in $\Omega$ such that $A \cap U$ is a complex submanifold of $\Omega$. A point $a \in A$ is called a singular point of $A$ if it is not regular.

Definition 2.2. Let $Y$ be a manifold and let $X, Z$ be submanifolds of $Y$. We say that the submanifolds $X$ and $Z$ intersect transversely if for every $x \in X \cap Z, T_{x}(X)+T_{x}(Z)=T_{x}(Y)$.
Assumption 2.3. Let $\tilde{M}$ be an analytic subset in an open neighborhood of the closed ball $\overline{\mathbf{B}}$. Furthermore, $\tilde{M}$ satisfies the following conditions:
(1) $\tilde{M}$ intersects $\partial \mathbf{B}$ transversely.
(2) $\tilde{M}$ has no singular points on $\partial \mathbf{B}$.
(3) $\tilde{M}$ is of pure dimension $d$, where $1 \leq d \leq n-1$.

Note that condition (3) implies that $\tilde{M}$ has no isolated singularities in B.

We emphasize that Assumption 2.3 will always be in force for the rest of the paper. Given such an $\tilde{M}$, we fix $M, \mathcal{R}, \mathcal{Q}$ and $Q$ as follows.

Notation 2.4. (a) Let $M=\tilde{M} \cap \mathbf{B}$.
(b) Denote $\mathcal{R}=\left\{f \in H_{n}^{2}: f=0\right.$ on $\left.M\right\}$.
(c) Denote $\mathcal{Q}=H_{n}^{2} \ominus \mathcal{R}$.
(d) Let $Q$ be the orthogonal projection from $H_{n}^{2}$ onto $\mathcal{Q}$.

In addition, we will simplify the notation for module operators used in the Introduction. Namely, we will write $Q_{\zeta_{j}}=\mathcal{Z}_{\mathcal{Q}, j}$ for $j=1, \ldots, n$. That is,

$$
Q_{\zeta_{j}}=Q M_{\zeta_{j}} \mid \mathcal{Q}
$$

$j=1, \ldots, n$. Thus the goal of the paper is to show that under the condition $d=1$, we have $\left[Q_{\zeta_{i}}, Q_{\zeta_{j}}^{*}\right] \in \mathcal{C}_{1}$ for all $i, j \in\{1, \ldots, n\}$. But this will take a very long journey.

Denote $S=\left\{z \in \mathbf{C}^{n}:|z|=1\right\}$, the unit sphere in $\mathbf{C}^{n}$. For $z \in \mathbf{C}^{n}$ and $r>0$, denote

$$
B(z, r)=\left\{w \in \mathbf{C}^{n}:|z-w|<r\right\} .
$$

By Assumption 2.3, there is an $s \in(0,1)$ such that

$$
\begin{equation*}
\mathcal{M}=\{z \in \tilde{M}: 1-s<|z|<1+s\} \tag{2.1}
\end{equation*}
$$

is a complex manifold of complex dimension $d$ and of finite volume.
For each $z \in \mathcal{M}$, let $T_{z}$ be the tangent space to $\mathcal{M}$ at the point $z$, viewed as a natural subspace of $\mathbf{C}^{n}$. For each $z \in \mathcal{M}$, let $p_{z}$ be the orthogonal projection of $z$ on $T_{z}$. Condition (1) in Assumption 2.3 says that if $z \in \tilde{M} \cap S$, then $p_{z} \neq 0$. Thus, reducing the value of $s \in(0,1)$ if necessary, we may assume that there is a $\gamma>0$ such that

$$
\begin{equation*}
\left|p_{z}\right| \geq \gamma \quad \text { for every } \quad z \in \mathcal{M} \tag{2.2}
\end{equation*}
$$

Denote

$$
\begin{equation*}
K=\{z \in \tilde{M}: 1-(s / 2) \leq|z| \leq 1\} \tag{2.3}
\end{equation*}
$$

Then $K$ is a compact subset of the complex manifold $\mathcal{M}$. By the standard facts known about such a pair of $\mathcal{M}$ and $K$ [23-25], the statements below hold true with constants that are independent of $z \in K$.

There are $a>0$ and $b>0$ such that for each $z \in K$, there is a map

$$
\begin{equation*}
G_{z}: T_{z} \cap B(0, a) \rightarrow \mathcal{M} \tag{2.4}
\end{equation*}
$$

that biholomorphically maps $T_{z} \cap B(0, a)$ onto an open subset of $\mathcal{M}$ with the properties that $G_{z}(0)=z$ and that

$$
\begin{equation*}
\left\{G_{z}(w): w \in T_{z} \cap B(0, a)\right\} \supset \mathcal{M} \cap B(z, b) \tag{2.5}
\end{equation*}
$$

Let $D G_{z}$ be the complex derivative of $G_{z}$. For each $w \in T_{z} \cap B(0, a)$, we have the local Taylor expansion

$$
G_{z}(w+u)=G_{z}(w)+\left(D G_{z}\right)(w) u+\int_{0}^{1}\left\{\left(D G_{z}\right)(w+t u)-\left(D G_{z}\right)(w)\right\} u d t
$$

$w+u \in T_{z} \cap B(0, a)$. In particular, at the point $w=0$ we have

$$
T_{z}=\left(D G_{z}\right)(0) T_{z}
$$

and

$$
G_{z}(u)=z+\left(D G_{z}\right)(0) u+\int_{0}^{1}\left\{\left(D G_{z}\right)(t u)-\left(D G_{z}\right)(0)\right\} u d t \quad \text { for } \quad u \in T_{z} \cap B(0, a)
$$

Reducing the values of $a$ and $b$ if necessary, we may assume that there are constants $0<\alpha \leq \beta<\infty$ such that for $w \in T_{z} \cap B(0, a)$, the linear transformation inequality

$$
\alpha \leq\left(D G_{z}\right)^{*}(w)\left(D G_{z}\right)(w) \leq \beta
$$

holds on $T_{z}$.
For each $z \in K$,

$$
T_{z}^{\perp}=\left\{u \in T_{z}:\left\langle u, p_{z}\right\rangle=0\right\}
$$

is a linear subspace of $T_{z}$ of dimension $d-1$. As a subspace of $\mathbf{C}^{n}, T_{z}^{\perp}$ is orthogonal to $z$.
Definition 2.5. (a) For each $z \in K$, we define

$$
T_{z}^{\bmod }=T_{z}^{\perp} \oplus\{\xi z: \xi \in \mathbf{C}\}
$$

which we consider as the modified complex tangent space at $z$.
(b) For each $z \in K$, let $P_{z}$ be the orthogonal projection from $\mathbf{C}^{n}$ onto $T_{z}^{\bmod }$.

Lemma 2.6. [28, Lemma 2.7] There exist $b_{0}>0$ and $c_{0}>0$ such that for every $z \in K$, $P_{z}$ is a biholomorphic map from $\mathcal{M} \cap B\left(z, b_{0}\right)$ onto an open set in $T_{z}^{\bmod }$ that contains $T_{z}^{\bmod } \cap B\left(z, c_{0}\right)$.

For $z \in K$, let $I_{z}: T_{z}^{\bmod } \cap B\left(z, c_{0}\right) \rightarrow \mathcal{M}$ be the inverse of $P_{z}$. For $x \in T_{z}^{\bmod } \cap B\left(z, c_{0}\right)$, the relation $P_{z} I_{z}(x)=x$ leads to

$$
\begin{equation*}
I_{z}(x)=x+h_{z}(x), \quad \text { where } \quad h_{z}(x)=I_{z}(x)-P_{z} I_{z}(x) \tag{2.6}
\end{equation*}
$$

That is, for each $z \in K, h_{z}$ maps $T_{z}^{\bmod } \cap B\left(z, c_{0}\right)$ into $\mathbf{C}^{n} \ominus T_{z}^{\bmod }$. We now fix a $c_{1} \in\left(0, c_{0}\right)$. By the analysis on page 8 in [28], there are constants $0<\alpha\left(c_{1}\right) \leq \beta\left(c_{1}\right)<\infty$ such that the operator inequality

$$
\begin{equation*}
\alpha\left(c_{1}\right) \leq\left(D I_{z}\right)^{*}(x)\left(D I_{z}\right)(x) \leq \beta\left(c_{1}\right) \tag{2.7}
\end{equation*}
$$

holds on the linear space $T_{z}^{\bmod }$ for all $z \in K$ and $x \in T_{z}^{\bmod } \cap B\left(z, c_{1}\right)$. Applying the standard open mapping theorem, there is a $0<b_{1}<b_{0}$ such that

$$
\begin{equation*}
\left\{I_{z}(x): x \in T_{z}^{\bmod } \cap B\left(z, c_{1}\right)\right\} \supset \mathcal{M} \cap B\left(z, b_{1}\right) \tag{2.8}
\end{equation*}
$$

Our analysis also involves the Bergman-metric structure of the ball. As usual, we write $\beta$ for the Bergman metric on $\mathbf{B}$. That is,

$$
\beta(z, w)=\frac{1}{2} \log \frac{1+\left|\varphi_{z}(w)\right|}{1-\left|\varphi_{z}(w)\right|}, \quad z, w \in \mathbf{B} .
$$

We recall that the Möbius transform $\varphi_{z}$ is given by the formula

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

when $z \neq 0$, and $\varphi_{0}(w)=-w$. For each $z \in \mathbf{B}$ and each $a>0$, we define the corresponding $\beta$-ball $D(z, a)=\{w \in \mathbf{B}: \beta(z, w)<a\}$.
Lemma 2.7. [28, Lemma 2.9] (1) Let $r>0$ be given. For each $\epsilon>0$, there is a $\delta=\delta(r, \epsilon)$ $\in(0,1)$ such that if $z \in K$ satisfies the condition $1-\delta \leq|z|<1$, then the inequality

$$
\beta\left(w, P_{z} w\right) \leq \epsilon
$$

holds for every $w \in D(z, r) \cap \mathcal{M}$.
(2) Let $z \in M \cap K$ and $r>0$ be such that $D(z, r / 2) \subset B\left(z, c_{0}\right)$ and $\beta\left(w, P_{z} w\right) \leq r / 3$ for every $w \in D(z, 2 r) \cap M$. Then $I_{z}\left(D(z, r / 2) \cap T_{z}^{\bmod }\right) \subset D(z, r) \cap M$.

For every $z \in K, T_{z}^{\bmod }$ is a $d$-dimensional linear subspace of $\mathbf{C}^{n}$. For convenience we will write $v$ for the natural volume measure on $T_{z}^{\text {mod }}$, even though for different $z \in K$ this may be a different linear subspace of $\mathbf{C}^{n}$. But since volume depends only on the Euclidean metric, which $T_{z}^{\text {mod }}$ inherits from $\mathbf{C}^{n}$, such a simplification of notation is justified.

For each $z \in K$, we have the Jacobian

$$
\begin{equation*}
J_{z}(x)=\operatorname{det}\left\{\left(D I_{z}\right)^{*}(x)\left(D I_{z}\right)(x)\right\} \tag{2.9}
\end{equation*}
$$

$x \in T_{z}^{\bmod } \cap B\left(z, c_{1}\right)$. Let $v_{\mathcal{M}}$ denote the natural volume measure on $\mathcal{M}$. Suppose that $z \in K$ and $U$ is an open set in $\mathcal{M} \cap B\left(z, b_{1}\right)$. By (2.8), we have $P_{z} U \subset T_{z}^{\bmod } \cap B\left(z, c_{1}\right)$. For any positive, continuous function $f$ on $U$, we have

$$
\begin{equation*}
\int_{U} f(w) d v_{\mathcal{M}}(w)=\int_{P_{z} U} f\left(I_{z}(x)\right) J_{z}(x) d v(x) \tag{2.10}
\end{equation*}
$$

Recall that this is in fact how volume is defined on $\mathcal{M}$.
In addition to the volume measure $v_{\mathcal{M}}$ on $\mathcal{M}$, we define the measure $v_{M}$ on $M=\tilde{M} \cap \mathbf{B}$ by the formula $v_{M}(E)=v_{\mathcal{M}}(E \cap \mathcal{M})$ for Borel sets $E \subset M$.

Lemma 2.8. [28, Lemma 2.10] Given any $a>0$ and $\kappa>-1$, there is a $0<C_{2.8}<\infty$ such that

$$
\int_{M} \frac{\left(1-|z|^{2}\right)^{a}\left(1-|w|^{2}\right)^{\kappa}}{|1-\langle w, z\rangle|^{d+1+a+\kappa}} d v_{M}(w) \leq C_{2.8}
$$

for every $z \in M$.
Moreover, it is known that if $\kappa>-1$, then

$$
\int_{M}\left(1-|w|^{2}\right)^{\kappa} d v_{M}(w)<\infty
$$

[28, page 15]. This finiteness is due to the fact that we can use the function $\rho(w)=1-|w|^{2}$ as one of the $2 d$ real coordinates on $M$ for $w \in M$ near $S$.

Lemma 2.9. [28, Lemma 2.11] Given any $a>0$ and $\kappa>-1$, there are $\delta>0$ and $0<C_{2.9}(\delta)<\infty$ such that

$$
\int_{M \backslash D(z, r)} \frac{\left(1-|z|^{2}\right)^{a}\left(1-|w|^{2}\right)^{\kappa}}{|1-\langle w, z\rangle|^{d+1+a+\kappa}} d v_{M}(w) \leq C_{2.9}(\delta) e^{-2 \delta r}
$$

for all $z \in M$ and $r>0$.
Following [28], we use a subscript $d$ to indicate a set in $\mathbf{C}^{d}$. For example, $\mathbf{B}_{d}=\{w \in$ $\left.\mathbf{C}^{d}:|w|<1\right\}$ and $D_{d}(z, r)=\left\{w \in \mathbf{B}_{d}: \beta(z, w)<r\right\}$. Similarly, $d v_{d}$ denotes the volume measure on $\mathbf{C}^{d}$. In particular, $d v_{1}$ is just the area measure on $\mathbf{C}$.

Lemma 2.10. Let $0<r<\infty$. If $f$ is a bounded analytic function on $D_{1}(z, r), z \in \mathbf{B}_{1}$, then

$$
\begin{equation*}
\int_{D_{1}(z, r)} f(w) \frac{1-|w|^{2}}{(1-\langle z, w\rangle)^{3}} d v_{1}(w)=\Phi(r) f(z) \tag{2.11}
\end{equation*}
$$

where

$$
\Phi(r)=\int_{D_{1}(0, r)}\left(1-|\zeta|^{2}\right) d v_{1}(\zeta)
$$

Proof. Let $w=\varphi_{z}(\zeta)$. By the formulas from [26, Theorem 2.2.2], we have

$$
1-\left\langle z, \varphi_{z}(\zeta)\right\rangle=\frac{1-|z|^{2}}{1-\langle z, \zeta\rangle} \quad \text { and } \quad 1-\left|\varphi_{z}(\zeta)\right|^{2}=\frac{\left(1-|z|^{2}\right)\left(1-|\zeta|^{2}\right)}{|1-\langle z, \zeta\rangle|^{2}} .
$$

Therefore the left-hand side of (2.11) equals

$$
\int_{D_{1}(0, r)} f\left(\varphi_{z}(\zeta)\right) \frac{\left(1-|z|^{2}\right)\left(1-|\zeta|^{2}\right)}{|1-\langle z, \zeta\rangle|^{2}}\left(\frac{1-\langle z, \zeta\rangle}{1-|z|^{2}}\right)^{3} \frac{\left(1-|z|^{2}\right)^{2}}{|1-\langle z, \zeta\rangle|^{4}} d v_{1}(\zeta)
$$

After the obvious cancellation, we find that

$$
\int_{D_{1}(z, r)} f(w) \frac{1-|w|^{2}}{(1-\langle z, w\rangle)^{3}} d v_{1}(w)=\int_{D_{1}(0, r)} \frac{f\left(\varphi_{z}(\zeta)\right)}{(1-\langle\zeta, z\rangle)^{3}}\left(1-|\zeta|^{2}\right) d v_{1}(\zeta)
$$

With respect to the Euclidean metric on $\mathbf{C}, D_{1}(0, r)$ is a disc centered at 0 . Hence the above equals $\Phi(r) f\left(\varphi_{z}(0)\right)(1-\langle 0, z\rangle)^{-3}=\Phi(r) f(z)$.
Lemma 2.11. [28, Lemma 3.2] For each given $0<r<\infty$, we have

$$
\lim _{t \uparrow 1} \sup \left\{\left|1-\frac{1-|x|^{2}}{1-\left|I_{z}(x)\right|^{2}}\right|:|z| \geq t, z \in M \text { and } x \in D(z, r) \cap T_{z}^{\bmod }\right\}=0
$$

and

$$
\limsup _{t \uparrow 1}\left\{\left|J_{z}(z)-J_{z}(x)\right|:|z| \geq t, z \in M \text { and } x \in D(z, r) \cap T_{z}^{\bmod }\right\}=0
$$

Lemma 2.12. Let $-1<\tau<\infty$. Then

$$
d \Omega(w)=\left(1-|w|^{2}\right)^{n-d+\tau} d v_{M}(w)
$$

is a Carleson measure for the weighted Bergman space $L_{a, \tau}^{2}=L_{a}^{2}\left(\mathbf{B},\left(1-|z|^{2}\right)^{\tau} d v(z)\right)$.
Proof. For each pair of $\zeta \in S$ and $r>0$, define $Q(\zeta, r)=\{z \in \mathbf{B}:|1-\langle z, \zeta\rangle|<r\}$. By the well-known [7, Theorem 1], to show that $\Omega$ is a Carleson measure for $L_{a, \tau}^{2}$, it suffices to find a $C_{1}$ such that

$$
\begin{equation*}
\Omega(Q(\zeta, r)) \leq C_{1} r^{n+1+\tau} \tag{2.12}
\end{equation*}
$$

for all $\zeta \in S$ and $r>0$. To prove this, note that

$$
\left(1-|w|^{2}\right)^{n-d+\tau}=\left(1-|w|^{2}\right)^{1+\tau} \cdot\left(1-|w|^{2}\right)^{n-1-d}
$$

If $w \in Q(\zeta, r) \cap M$, then $1-|w|^{2} \leq 2 r$. Since $1+\tau>0$, we have
$\Omega(Q(\zeta, r))=\int_{Q(\zeta, r) \cap M}\left(1-|w|^{2}\right)^{n-d+\tau} d v_{M}(w) \leq(2 r)^{1+\tau} \int_{Q(\zeta, r) \cap M}\left(1-|w|^{2}\right)^{n-1-d} d v_{M}(w)$.
By inequality (2.26) in [28],

$$
\int_{Q(\zeta, r) \cap M}\left(1-|w|^{2}\right)^{n-1-d} d v_{M}(w) \leq C r^{n}
$$

for all $\zeta \in S$ and $r>0$. Thus (2.12) indeed holds.
Under the usual identification of $\mathbf{C}$ with $\mathbf{R}^{2}$, we can also view $T_{z}$ as a subspace of $\mathbf{R}^{2 n}$ of real dimension $2 d$, equipped with the real inner product. Thus if $z \in \mathcal{M}$ and $h$ is a real-valued $C^{1}$-function on an open neighborhood $U$ of $z$ in $\mathbf{C}^{n} \cong \mathbf{R}^{2 n}$, then we define
$\left(\nabla_{\mathcal{M}} h\right)(z)$ to be the orthogonal projection of the real vector $(\nabla h)(z)$ onto the real subspace $T_{z}$. If $h$ is complex-valued, we can write $h=h_{1}+i h_{2}$, where $h_{1}$ and $h_{2}$ are real valued. In this case, we define $\left(\nabla_{\mathcal{M}} h\right)(z)=\left(\nabla_{\mathcal{M}} h_{1}\right)(z)+i\left(\nabla_{\mathcal{M}} h_{2}\right)(z)$. This defines the operation $\nabla_{\mathcal{M}}$. We think of $\nabla_{\mathcal{M}}$ as the gradient in the directions tangent to $\mathcal{M}$.

For each $0<t \leq 1$, we define the sets

$$
\begin{equation*}
M^{(t)}=\left\{z \in M: 1-|z|^{2}<t\right\} \quad \text { and } \quad N^{(t)}=\left\{z \in M: 1-|z|^{2}>t\right\} \tag{2.13}
\end{equation*}
$$

which will appear frequently in the sequel. For each $-1<\tau<\infty$, we denote

$$
d \nu_{\tau}(w)=\left(1-|w|^{2}\right)^{\tau} d v_{M}(w) .
$$

Lemma 2.13. Given any $-1<\tau<\infty$ and $0<t \leq 1$, there are constants $0<a<b<t$ and $0<C<\infty$ such that the inequality

$$
\begin{aligned}
\int_{M^{(t)}}|f(w)|^{2} d \nu_{\tau}(w) & \leq C \int_{M^{(b)}}\left|\left(\nabla_{\mathcal{M}} f\right)(w)\right|^{2}\left(1-|w|^{2}\right)^{2} d \nu_{\tau}(w) \\
& +C \int_{N^{(a)} \cap M^{(t)}}|f(w)|^{2} d \nu_{\tau}(w)
\end{aligned}
$$

holds for every $C^{1}$ function $f$ on any open set containing the closure of $M^{(t)}$.
The proof of this lemma is essentially the same as the proof of [29, Lemma 3.1]. For that reason we leave the proof of Lemma 2.13 to Appendix 1.

As usual, we write $\partial=\left(\partial_{1}, \ldots, \partial_{n}\right)$, the analytic gradient on $\mathbf{C}^{n}$. By the multi-index convention that we follow, for $\alpha=\left(\alpha_{1}, \ldots, \alpha_{n}\right) \in \mathbf{Z}_{+}^{n}, \partial^{\alpha}$ denotes $\partial_{1}^{\alpha_{1}} \cdots \partial_{n}^{\alpha_{n}}$.

Lemma 2.14. Let $-1<\tau<\infty$ and $k \in \mathbf{N}$ be given. There are constants $0<a<b<1$ and $0<C<\infty$ such that if $f$ is any analytic function on an open set containing the closure of $M$, then

$$
\begin{aligned}
\int_{M}|f(w)|^{2} d \nu_{\tau}(w) & \leq C \sum_{|\alpha|=k} \int_{M^{(b)}}\left|\left(\partial^{\alpha} f\right)(w)\right|^{2}\left(1-|w|^{2}\right)^{2 k} d \nu_{\tau}(w) \\
& +C \sum_{0 \leq|\beta| \leq k-1} \int_{N^{(a)}}\left|\left(\partial^{\beta} f\right)(w)\right|^{2} d \nu_{\tau}(w)
\end{aligned}
$$

Proof. This follows from Lemma 2.13 by an obvious induction on $k$.
Let us recall the family of spaces $\mathcal{H}^{(t)}$ introduced in [19]. For each $-n \leq t<\infty$, let $\mathcal{H}^{(t)}$ be the Hilbert space of analytic functions on $\mathbf{B}$ which has the function

$$
K_{w}^{(t)}(z)=\frac{1}{(1-\langle z, w\rangle)^{n+1+t}}, \quad z, w \in \mathbf{B}
$$

as its reproducing kernel. Equivalently, $\mathcal{H}^{(t)}$ is the completion of $\mathbf{C}\left[z_{1}, \ldots, z_{n}\right]$ with respect to the norm $\|\cdot\|_{t}$ arising from the inner product $\langle\cdot, \cdot\rangle_{t}$ defined according to the following rules: $\left\langle z^{\alpha}, z^{\beta}\right\rangle_{t}=0$ whenever $\alpha \neq \beta$,

$$
\begin{equation*}
\left\langle z^{\alpha}, z^{\alpha}\right\rangle_{t}=\frac{\alpha!}{\prod_{j=1}^{|\alpha|}(n+t+j)} \tag{2.14}
\end{equation*}
$$

if $\alpha \in \mathbf{Z}_{+}^{n} \backslash\{0\}$, and $\langle 1,1\rangle_{t}=1$. It is well known that $H_{n}^{2}=\mathcal{H}^{(-n)}$ and that for each $-1<t<\infty, \mathcal{H}^{(t)}$ is a weighted Bergman space on $\mathbf{B}$.

Recall that the formula

$$
R=z_{1} \partial_{1}+\cdots+z_{n} \partial_{n}
$$

defines the radial derivative on $\mathbf{B}$. Let $m \in \mathbf{N}$ and $-n \leq t<\infty$ satisfy the condition $2 m+t>-1$. For such a pair of $m$ and $t$, we define

$$
\|f\|_{m, t}^{2}=|f(0)|^{2}+\int_{\mathbf{B}}\left|\left(R^{m} f\right)(z)\right|^{2}\left(1-|z|^{2}\right)^{2 m+t} d v(z)
$$

whenever $f$ is an analytic function on $\mathbf{B}$. The following is well known:
Lemma 2.15. Let $m \in \mathbf{N}$ and $-n \leq t<\infty$ satisfy the condition $2 m+t>-1$. Then there exist constants $0<c \leq C<\infty$ such that

$$
c\|f\|_{m, t} \leq\|f\|_{t} \leq C\|f\|_{m, t}
$$

for every analytic function $f$ on $\mathbf{B}$.
Proof. When $\alpha \neq \beta$, we have $\left\langle z^{\alpha}, z^{\beta}\right\rangle_{t}=0$ by definition. When $\alpha \neq \beta$, it is easy to see that $\left\langle z^{\alpha}, z^{\beta}\right\rangle_{m, t}=0$, where $\langle\cdot, \cdot\rangle_{m, t}$ is the inner product that corresponds to the norm $\|\cdot\|_{m, t}$ defined above. Thus it suffices to find constants $0<c \leq C<\infty$ such that

$$
\begin{equation*}
c\left\|z^{\alpha}\right\|_{m, t} \leq\left\|z^{\alpha}\right\|_{t} \leq C\left\|z^{\alpha}\right\|_{m, t} \tag{2.15}
\end{equation*}
$$

for every $\alpha \in \mathbf{Z}_{+}^{n}$. We have $R^{m} z^{\alpha}=|\alpha|^{m} z^{\alpha}$. Therefore, for any $\alpha \neq 0$,

$$
\begin{align*}
\left\|z^{\alpha}\right\|_{m, t}^{2} & =|\alpha|^{2 m} \int_{\mathbf{B}}\left|z^{\alpha}\right|^{2}\left(1-|z|^{2}\right)^{2 m+t} d v(z) \\
& =\frac{|\alpha|^{2 m}(n-1)!\alpha!}{(n-1+|\alpha|)!} 2 n \int_{0}^{1} r^{2|\alpha|+2 n-1}\left(1-r^{2}\right)^{2 m+t} d r \\
& =\frac{|\alpha|^{2 m} n!\alpha!}{(n-1+|\alpha|)!} \int_{0}^{1} x^{|\alpha|+n-1}(1-x)^{2 m+t} d x \\
& =\frac{|\alpha|^{2 m} n!\alpha!}{(n-1+|\alpha|)!} \cdot \frac{(|\alpha|+n-1)!}{\prod_{j=1}^{|\alpha|+n}(2 m+t+j)}=\frac{|\alpha|^{2 m} n!\alpha!}{\prod_{j=1}^{|\alpha|+n}(2 m+t+j)} . \tag{2.16}
\end{align*}
$$

By Stirling's asymptotic formula (see, e.g., identity (3.3) in [18]), we have

$$
\begin{aligned}
& \prod_{j=1}^{|\alpha|+n}(2 m+t+j) \approx(2 m+t+|\alpha|+n)^{2 m+t+|\alpha|+n+(1 / 2)} e^{-|\alpha|} \quad \text { whereas } \\
& \quad \prod_{j=1}^{|\alpha|}(n+t+j) \approx(n+t+|\alpha|)^{n+t+|\alpha|+(1 / 2)} e^{-|\alpha|}
\end{aligned}
$$

Combining these formulas with (2.14) and (2.16), we obtain (2.15).
Lemma 2.16. Given $m \in \mathbf{N}$ and $t>-1$, there is a constant $0<C<\infty$ such that

$$
\begin{equation*}
\int_{\mathbf{B}}\left|\left(\partial^{\beta} f\right)(z)\right|^{2}\left(1-|z|^{2}\right)^{t} d v(z) \leq C \int_{\mathbf{B}}\left|\left(R^{m} f\right)(z)\right|^{2}\left(1-|z|^{2}\right)^{t} d v(z) \tag{2.17}
\end{equation*}
$$

for every $\beta \in \mathbf{Z}_{+}^{n}$ satisfying the condition $|\beta|=m$ and every analytic function $f$ on $\mathbf{B}$.
Proof. Similar to what happened in the previous proof, it suffices to find a $0<C<\infty$ such that

$$
\begin{equation*}
\int_{\mathbf{B}}\left|\partial^{\beta} z^{\alpha}\right|^{2}\left(1-|z|^{2}\right)^{t} d v(z) \leq C \int_{\mathbf{B}}\left|R^{m} z^{\alpha}\right|^{2}\left(1-|z|^{2}\right)^{t} d v(z) \tag{2.18}
\end{equation*}
$$

for every $\alpha \in \mathbf{Z}_{+}^{n}$. Write $\beta=\left(\beta_{1}, \ldots, \beta_{n}\right)$. Note that for any $\alpha=\left(\alpha_{1}, \ldots, \alpha_{n}\right) \in \mathbf{Z}_{+}^{n}$, the left-hand side of (2.18) is 0 unless $\alpha_{j} \geq \beta_{j}$ for every $j \in\{1, \ldots, n\}$. Suppose that this condition is satisfied. Then $\alpha-\beta \in \mathbf{Z}_{+}^{n}$ and we have

$$
\partial^{\beta} z^{\alpha}=z^{\alpha-\beta} \prod_{\beta_{\nu}>0} \prod_{i_{\nu}=1}^{\beta_{\nu}}\left(\alpha_{\nu}-\beta_{\nu}+i_{\nu}\right) .
$$

Therefore

$$
\begin{align*}
\int_{\mathbf{B}}\left|\partial^{\beta} z^{\alpha}\right|^{2}\left(1-|z|^{2}\right)^{t} d v(z) & =\prod_{\beta_{\nu}>0} \prod_{i_{\nu}=1}^{\beta_{\nu}}\left(\alpha_{\nu}-\beta_{\nu}+i_{\nu}\right)^{2} \int_{\mathbf{B}}\left|z^{\alpha-\beta}\right|^{2}\left(1-|z|^{2}\right)^{t} d v(z) \\
& =\prod_{\beta_{\nu}>0} \prod_{i_{\nu}=1}^{\beta_{\nu}}\left(\alpha_{\nu}-\beta_{\nu}+i_{\nu}\right)^{2} \cdot \frac{n!(\alpha-\beta)!}{\prod_{j=1}^{|\alpha-\beta|+n}(t+j)} \\
& \leq \frac{n!\alpha!|\alpha|^{m} \prod_{j=|\alpha-\beta|+n+1}^{|\alpha|+n}(t+j)}{\prod_{j=1}^{|\alpha|+n}(t+j)} . \tag{2.19}
\end{align*}
$$

On the other hand, since $R^{m} z^{\alpha}=|\alpha|^{m} z^{\alpha}$, we have

$$
\begin{equation*}
\int_{\mathbf{B}}\left|R^{m} z^{\alpha}\right|^{2}\left(1-|z|^{2}\right)^{t} d v(z)=\frac{n!\alpha!|\alpha|^{2 m}}{\prod_{j=1}^{|\alpha|+n}(t+j)} \tag{2.20}
\end{equation*}
$$

Since $|\beta|=m$, we have $|\alpha|=|\alpha-\beta|+m$, and consequently $\prod_{j=|\alpha-\beta|+n+1}^{|\alpha|+n}(t+j) \leq C|\alpha|^{m}$. Combining this inequality with (2.19) and (2.20), we obtain (2.18).

We end the section with one more notation. For any $u=\left(u_{1}, \ldots, u_{n}\right) \in \mathbf{C}^{n}$, we write

$$
\begin{equation*}
\partial_{u}=u_{1} \partial_{1}+\cdots+u_{n} \partial_{n} \quad \text { and } \quad \overline{\partial_{u}}=\bar{u}_{1} \bar{\partial}_{1}+\cdots+\bar{u}_{n} \bar{\partial}_{n} \tag{2.21}
\end{equation*}
$$

## 3. Spectral gap

For the rest of the paper, we assume $d=1$. Accordingly, we define the measure

$$
\begin{equation*}
d \mu(w)=\left(1-|w|^{2}\right) d v_{M}(w) . \tag{3.1}
\end{equation*}
$$

We fix a $t_{0} \in(0,1)$ such that

$$
\begin{equation*}
M^{\left(2 t_{0}\right)} \subset K \tag{3.2}
\end{equation*}
$$

(see (2.3) and (2.13)). Recall that for $z \in \mathcal{M}, p_{z}$ is the orthogonal projection of $z$ on $T_{z}$.
Definition 3.1. For any $C^{1}$ function $f$ on an open set containing $M$, we define

$$
\|f\|_{*}=\left\{\int_{M^{\left(t_{0}\right)}}\left|\left(\partial_{p_{w}} f\right)(w)\right|^{2} d \mu(w)\right\}^{1 / 2}
$$

Proposition 3.2. There is a $0<C<\infty$ such that

$$
\begin{equation*}
\int_{M}\left|\left(\partial_{j} f\right)(w)\right|^{2} d \mu(w) \leq C\|f\|^{2} \tag{3.3}
\end{equation*}
$$

for all $f \in H_{n}^{2}$ and $1 \leq j \leq n$.
Proof. We pick a natural number $m \geq 2$ such that $2 m-n \geq 0$. Given an $f \in H_{n}^{2}$, define the function $f_{r}(w)=f(r w)$ for each $0<r<1$. Each $f_{r}$ is analytic on an open ball containing $\overline{\mathbf{B}}$. By Lemma 2.14, we have

$$
\begin{aligned}
\int_{M}\left|\left(\partial_{j} f_{r}\right)(w)\right|^{2} d \mu(w) & \leq C_{1} \sum_{|\alpha|=m-1} \int_{M^{(b)}}\left|\left(\partial^{\alpha} \partial_{j} f_{r}\right)(w)\right|^{2}\left(1-|w|^{2}\right)^{2 m-2} d \mu(w) \\
& +C_{1} \sum_{0 \leq|\beta| \leq m-2} \int_{N^{(a)}}\left|\left(\partial^{\beta} \partial_{j} f_{r}\right)(w)\right|^{2} d \mu(w)
\end{aligned}
$$

By the definition of $N^{(a)}$, the second term is dominated by $C_{2}\left\|f_{r}\right\|^{2} \leq C_{2}\|f\|^{2}$. It follows from Lemma 2.12 that for each $\alpha \in \mathbf{Z}_{+}^{n}$ with $|\alpha|=m-1$,

$$
\begin{aligned}
& \int_{M^{(b)}}\left|\left(\partial^{\alpha} \partial_{j} f_{r}\right)(w)\right|^{2}\left(1-|w|^{2}\right)^{2 m-2} d \mu(w) \\
& \quad \leq C_{3} \int_{\mathbf{B}}\left|\left(\partial^{\alpha} \partial_{j} f_{r}\right)(w)\right|^{2}\left(1-|w|^{2}\right)^{2 m-n} d v(w) \leq C_{4}\left\|f_{r}\right\|^{2} \leq C_{4}\|f\|^{2}
\end{aligned}
$$

where the second $\leq$ follows from Lemmas 2.15 and 2.16. Thus

$$
\int_{M}\left|\left(\partial_{j} f_{r}\right)(w)\right|^{2} d \mu(w) \leq C_{5}\|f\|^{2}
$$

Applying Fatou's lemma,

$$
\int_{M}\left|\left(\partial_{j} f\right)(w)\right|^{2} d \mu(w) \leq \liminf _{r \uparrow 1} \int_{M}\left|\left(\partial_{j} f_{r}\right)(w)\right|^{2} d \mu(w) \leq C_{5}\|f\|^{2}
$$

This completes the proof.
We will now try to modify the reproducing kernel for $H_{n}^{2}$ given by (1.1). For each pair of $w \in \mathbf{B}$ and $u \in \mathbf{C}^{n}$, we define

$$
\begin{equation*}
K_{w, u}(\zeta)=\frac{\langle\zeta, u\rangle}{(1-\langle\zeta, w\rangle)^{2}} \tag{3.4}
\end{equation*}
$$

For $w \in \mathbf{B}, u \in \mathbf{C}^{n}$ and $f \in H_{n}^{2}$, it is easy to see that

$$
\begin{equation*}
\left\langle f, K_{w, u}\right\rangle=\left.\frac{d}{d t}\left\langle f, K_{w+t u}\right\rangle\right|_{t=0}=\left.\frac{d}{d t} f(w+t u)\right|_{t=0}=\left(\partial_{u} f\right)(w) \tag{3.5}
\end{equation*}
$$

That is, $K_{w, u}$ is the reproducing kernel for the directional derivative $\partial_{u}$.
Lemma 3.3. Let $w \in M^{\left(t_{0}\right)}$. If $u \in T_{w}$, then $K_{w, u} \in \mathcal{Q}$.
Proof. Since $u \in T_{w}$, there is a smooth path $\gamma:(-c, c) \rightarrow M^{\left(t_{0}\right)}$ such that $\gamma(0)=w$ and $\gamma^{\prime}(0)=u$. Thus

$$
K_{w, u}=\left.\frac{d}{d t} K_{\gamma(t)}\right|_{t=0}
$$

Let $f \in \mathcal{R}$. Since the range of $\gamma$ is contained in $M$, we have $\left\langle f, K_{\gamma(t)}\right\rangle=0$ for every $t \in(-c, c)$. Therefore

$$
\left\langle f, K_{w, u}\right\rangle=\left.\frac{d}{d t}\left\langle f, K_{\gamma(t)}\right\rangle\right|_{t=0}=\left.\frac{d}{d t} 0\right|_{t=0}=0
$$

This shows that $K_{w, u} \perp \mathcal{R}$. That is, $K_{w, u} \in \mathcal{Q}$.
We now define the operator

$$
\begin{equation*}
T_{1}=\int_{M^{\left(t_{0}\right)}} K_{w, p_{w}} \otimes K_{w, p_{w}} d \mu(w) \tag{3.6}
\end{equation*}
$$

Proposition 3.4. The operator $T_{1}$ defined above is bounded on the Drury-Arveson space $H_{n}^{2}$. Moreover, $T_{1}$ maps $H_{n}^{2}$ into the quotient module $\mathcal{Q}$.
Proof. By (3.5), for $f \in H_{n}^{2}$ we have $\left\langle T_{1} f, f\right\rangle=\|f\|_{*}^{2}$. Hence the boundedness of $T_{1}$ follows from Proposition 3.2. For $w \in M^{\left(t_{0}\right)}$, since $p_{w} \in T_{w}$, Lemma 3.3 tells us that $K_{w, p_{w}} \in \mathcal{Q}$. Therefore $T_{1}$ maps $H_{n}^{2}$ into $\mathcal{Q}$.

Obviously, $T_{1}$ is a positive operator on $H_{n}^{2}$. The main goal of this section is to prove that there is an important gap in the spectrum of $T_{1}$ :

Theorem 3.5. There is a $c>0$ such that the spectrum of $T_{1}$ does not intersect $(0, c)$.
The proof of Theorem 3.5 requires some preparations.
Since we now assume $d=1, B_{1}(0, r)=\{\xi \in \mathbf{C}:|\xi|<r\}, r>0$. Since $\mathcal{M}$ is a complex manifold with $\operatorname{dim}_{\mathbf{C}} \mathcal{M}=1$, for each $y \in \mathcal{M}$, there is an open subset $V_{y}$ of $\mathcal{M}$ containing $y$ and a biholomorphic map

$$
\rho_{y}: B_{1}(0,2) \rightarrow V_{y}
$$

such that $\rho_{y}(0)=y$. For each $\xi \in B_{1}(0,2), \rho_{y}^{\prime}(\xi)$ is obviously a complex tangent vector to $\mathcal{M}$ at the point $\rho_{y}(\xi)$, and we have $\rho_{y}^{\prime}(\xi) \neq 0$ since $\rho_{y}$ is biholomorphic. Since $\operatorname{dim}_{\mathbf{C}}\left(T_{\rho_{y}(\xi)}\right)=1$, we have $T_{\rho_{y}(\xi)}=\mathbf{C} \rho_{y}^{\prime}(\xi)$ for every $\xi \in B_{1}(0,2)$. Define

$$
\eta_{y}(w)=\rho_{y}^{\prime}\left(\rho_{y}^{-1}(w)\right), \quad w \in V_{y} .
$$

Then $T_{w}=\mathbf{C} \eta_{y}(w)$ for every $w \in V_{y}$. The important fact to keep in mind is that $\eta_{y}$ is analytic on $V_{y}$. Since $p_{w}$ is the orthogonal projection of $w$ onto $T_{w}$, there is a continuous function $s_{y}: V_{y} \rightarrow \mathbf{C}$ such that

$$
\begin{equation*}
p_{w}=s_{y}(w) \eta_{y}(w) \quad \text { for every } \quad w \in V_{y} \tag{3.7}
\end{equation*}
$$

This identity embodies one of our main observations: modulo continuous scalar multiples, $p_{w}$ is locally analytic. This fact will be crucial for the proof of Theorem 3.5.

We now define $U_{y}=\rho_{y} B_{1}(0,1)$, which is an open subset of $\mathcal{M}$ containing $y$. Obviously, $\overline{U_{y}}$ is a compact subset of $V_{y}$. Thus $s_{y}$ is uniformly continuous on $\overline{U_{y}}$. By (2.2), the infimum of $\left|s_{y}\right|$ on the set $\overline{U_{y}}$ is greater than 0 . Since $K$ is a compact subset of $\mathcal{M}$, there is a finite subset $F$ of $K$ such that

$$
\bigcup_{y \in F} U_{y} \supset K
$$

By general topology, this implies
Lemma 3.6. There is an $\epsilon>0$ such that for each $\zeta \in K$, there is a $y=y(\zeta) \in F$ for which the containment $\{w \in \mathcal{M}:|\zeta-w|<\epsilon\} \subset U_{y}$ holds.
Lemma 3.7. For each $y \in M^{\left(t_{0}\right)}$, there is an open neighborhood $N_{y}$ of $y$ in $M^{\left(t_{0}\right)}$ which has the following property. Let $\left\{f_{k}\right\}$ be a sequence in $H_{n}^{2}$. If the sequence $\left\{T_{1}^{1 / 2} f_{k}\right\}$ weakly converges to 0 , then

$$
\lim _{k \rightarrow \infty} \sup \left\{\left|\left(\partial_{p_{w}} f_{k}\right)(w)\right|: w \in N_{y}\right\}=0
$$

Proof. For each $y \in M^{\left(t_{0}\right)}$, consider the biholomorphic map $\rho_{y}: B_{1}(0,2) \rightarrow V_{y}$ introduced above. Recall that $\rho_{y}(0)=y$. Since $y$ is now in $M^{\left(t_{0}\right)}$, there are $\epsilon=\epsilon(y)>0$ and $r=r(y) \in(0,1)$ such that $\rho_{y}\left(B_{1}(0, r)\right) \subset\left\{w \in M^{\left(t_{0}\right)}: 1-|w|^{2}>\epsilon\right\}$. We will show that the lemma holds for the open set $N_{y}=\rho_{y}\left(B_{1}(0, r / 2)\right)$.

We begin with the Bergman space $L_{a}^{2}\left(B_{1}(0, r), d v_{1}\right)$. For each $f \in H_{n}^{2}$, define

$$
(G f)(\xi)=\left(\partial_{\rho_{y}^{\prime}(\xi)} f\right)\left(\rho_{y}(\xi)\right), \quad \xi \in B_{1}(0, r)
$$

Using the conditions $r<1$ and (3.7), we have

$$
\begin{array}{rl}
\int_{B_{1}(0, r)}|(G f)(\xi)|^{2} & d v_{1}(\xi) \leq C_{1} \int_{B_{1}(0, r)}\left|s_{y}\left(\rho_{y}(\xi)\right)\right|^{2}\left|\left(\partial_{\rho_{y}^{\prime}(\xi)} f\right)\left(\rho_{y}(\xi)\right)\right|^{2}\left|\rho_{y}^{\prime}(\xi)\right|^{2} d v_{1}(\xi) \\
& =C_{1} \int_{\rho_{y}\left(B_{1}(0, r)\right)}\left|\left(\partial_{p_{w}} f\right)(w)\right|^{2} d v_{M}(w) \\
& \leq C_{1} \epsilon^{-1} \int_{\rho_{y}\left(B_{1}(0, r)\right)}\left|\left(\partial_{p_{w}} f\right)(w)\right|^{2}\left(1-|w|^{2}\right) d v_{M}(w) \\
& =C_{1} \epsilon^{-1} \int_{\rho_{y}\left(B_{1}(0, r)\right)}\left|\left(\partial_{p_{w}} f\right)(w)\right|^{2} d \mu(w) \\
& \leq C_{1} \epsilon^{-1}\left\langle T_{1} f, f\right\rangle=C_{1} \epsilon^{-1}\left\|T_{1}^{1 / 2} f\right\|^{2} \tag{3.8}
\end{array}
$$

Thus there is a bounded operator $W: H_{n}^{2} \rightarrow L_{a}^{2}\left(B_{1}(0, r), d v_{1}\right)$ such that $G=W T_{1}^{1 / 2}$.
Now let $\left\{f_{k}\right\}$ be any sequence in $H_{n}^{2}$ such that $\left\{T_{1}^{1 / 2} f_{k}\right\}$ weakly converges to 0 . Since $G=W T_{1}^{1 / 2}$, the sequence $\left\{G f_{k}\right\}$ weakly converges to 0 in $L_{a}^{2}\left(B_{1}(0, r), d v_{1}\right)$. Using the reproducing kernel for the Bergman space, we have

$$
\lim _{k \rightarrow \infty} \sup \left\{\left|\left(\partial_{\rho_{y}^{\prime}(\xi)} f_{k}\right)\left(\rho_{y}(\xi)\right)\right|: \xi \in B_{1}(0, r / 2)\right\}=0
$$

By (3.7) and the boundedness of $s_{y} \circ \rho_{y}$ on $B_{1}(0, r / 2)$, the above limit implies

$$
\lim _{k \rightarrow \infty} \sup \left\{\left|\left(\partial_{p_{w}} f_{k}\right)(w)\right|: w \in N_{y}\right\}=0
$$

as promised.
Lemma 3.8. Define the operators $B$ and $B_{r}$ on $L^{2}(M, d \mu)$ by the formulas

$$
\begin{aligned}
(B f)(z) & =\int_{M} \frac{f(w)}{|1-\langle z, w\rangle|^{3}} d \mu(w) \quad \text { and } \\
\left(B_{r} f\right)(z) & =\int_{M \backslash D(z, r)} \frac{f(w)}{|1-\langle z, w\rangle|^{3}} d \mu(w)
\end{aligned}
$$

for $f \in L^{2}(M, d \mu), r>0$. Then $\|B\|<\infty$ and $\left\|B_{r}\right\| \rightarrow 0$ as $r \rightarrow \infty$.
Proof. We set $a=1 / 2$ and $\kappa=1 / 2$. Define $h(w)=\left(1-|w|^{2}\right)^{-1 / 2}, w \in M$. Then

$$
\left(B_{r} h\right)(z)=\int_{M \backslash D(z, r)} \frac{\left(1-|w|^{2}\right)^{\kappa}}{|1-\langle z, w\rangle|^{1+1+a+\kappa}} d v_{M}(w) .
$$

By Lemma 2.9, we have $\left(B_{r} h\right)(z) \leq C_{2.9}(\delta) e^{-2 \delta r}\left(1-|z|^{2}\right)^{-a}=C_{2.9}(\delta) e^{-2 \delta r} h(z), z \in M$. Since the kernel function $|1-\langle z, w\rangle|^{-3}$ is symmetric with respect to $z$ and $w$, we can now apply the Schur test to conclude that $\left\|B_{r}\right\| \leq C_{2.9}(\delta) e^{-2 \delta r}$. Hence $\left\|B_{r}\right\| \rightarrow 0$ as $r \rightarrow \infty$.

Similarly, by Lemma 2.8 we have $(B h)(z) \leq C_{2.8} h(z), z \in M$. Thus it follows from the Schur test that $\|B\| \leq C_{2.8}$. This completes the proof.

For each $f \in H_{n}^{2}$, define

$$
(X f)(z)=\int_{M^{\left(t_{0}\right)}} \frac{\left(\partial_{p_{w}} f\right)(w)\left\langle p_{z}, w\right\rangle\left\langle z, p_{w}\right\rangle}{(1-\langle z, w\rangle)^{3}} d \mu(w), \quad z \in \mathbf{B} .
$$

Lemma 3.9. Given any $\delta>0$, there exist constants $0<\tau<t_{0}$ and $0<C<\infty$ such that

$$
\int_{M^{(t)}}\left|\left(\partial_{p_{z}} f\right)(z)\right|^{2} d \mu(z) \leq C \int_{M^{(t)}}|(X f)(z)|^{2} d \mu(z)+\delta\|f\|_{*}^{2}
$$

for all $0<t \leq \tau$ and $f \in H_{n}^{2}$.
Proof. We begin with a large $1 \leq r<\infty$, whose exact value will be determined below. With such an $r$, there is a $0<\tau_{1} \leq t_{0}$ such that if $0<t \leq \tau_{1}$, then for $z \in M^{(t)}$ we have $D(z, 2 r) \subset B\left(z, \min \left\{b_{1}, c_{1}\right\}\right)$ (see (2.8)). By Lemma 2.7(1), there is a $0<\tau_{2} \leq \tau_{1}$ such that if $0<t \leq \tau_{2}$, then for $z \in M^{(t)}$ and $w \in D(z, r) \cap M$ we have $\beta\left(w, P_{z} w\right)<r$. Thus $P_{z} w \in D(z, 2 r) \cap T_{z}^{\bmod }$ and $I_{z}\left(P_{z} w\right)=w \in D(z, r) \cap M$. That is, if $0<t \leq \tau_{2}$, then

$$
\begin{equation*}
I_{z}\left(D(z, 2 r) \cap T_{z}^{\bmod }\right) \supset D(z, r) \cap M \quad \text { for every } \quad z \in M^{(t)} \tag{3.9}
\end{equation*}
$$

By (2.7), there is a constant $1 \leq C_{1}<\infty$ such that the inequality

$$
\begin{equation*}
\left|I_{z}(x)-I_{z}\left(x^{\prime}\right)\right| \leq C_{1}\left|x-x^{\prime}\right| \tag{3.10}
\end{equation*}
$$

holds for every triple of $z \in K$ and $x, x^{\prime} \in T_{z}^{\bmod } \cap B\left(z, c_{1}\right)$. Therefore there is a $0<\tau_{3} \leq \tau_{2}$ such that

$$
\begin{equation*}
I_{z}\left(D(z, 2 r) \cap T_{z}^{\bmod }\right) \subset M^{\left(t_{0}\right)} \quad \text { if } \quad z \in M^{\left(\tau_{3}\right)} \tag{3.11}
\end{equation*}
$$

Let us write $U(z)=I_{z}\left(D(z, 2 r) \cap T_{z}^{\bmod }\right)$ for $z \in M^{\left(\tau_{3}\right)}$.
Let $f \in H_{n}^{2}$ be given. By (3.11), for $z \in M^{\left(\tau_{3}\right)}$ we have

$$
(X f)(z)=\left\langle p_{z}, z\right\rangle A(z)+B(z)+C(z),
$$

where

$$
\begin{aligned}
& A(z)=\int_{U(z)}\left(\partial_{p_{w}} f\right)(w)\left\langle z, p_{w}\right\rangle \frac{1-|w|^{2}}{(1-\langle z, w\rangle)^{3}} d v_{M}(w) \\
& B(z)=\int_{U(z)}\left(\partial_{p_{w}} f\right)(w)\left\langle p_{z}, w-z\right\rangle\left\langle z, p_{w}\right\rangle \frac{1-|w|^{2}}{(1-\langle z, w\rangle)^{3}} d v_{M}(w) \quad \text { and } \\
& C(z)=\int_{M^{\left(t_{0}\right)} \backslash U(z)}\left(\partial_{p_{w}} f\right)(w)\left\langle p_{z}, w\right\rangle\left\langle z, p_{w}\right\rangle \frac{1-|w|^{2}}{(1-\langle z, w\rangle)^{3}} d v_{M}(w)
\end{aligned}
$$

Since $P_{z} U(z)=D(z, 2 r) \cap T_{z}^{\bmod }, z \in M^{\left(\tau_{3}\right)}$, by (2.10) we have

$$
A(z)=\int_{D(z, 2 r) \cap T_{z}^{\bmod }}\left(\partial_{p_{I_{z}(x)}} f\right)\left(I_{z}(x)\right)\left\langle z, p_{I_{z}(x)}\right\rangle \frac{1-\left|I_{z}(x)\right|^{2}}{\left(1-\left\langle z, I_{z}(x)\right\rangle\right)^{3}} J_{z}(x) d v_{1}(x)
$$

Recall from (2.6) that $\left\langle z, I_{z}(x)\right\rangle=\langle z, x\rangle$. Writing

$$
F(z, x)=1-\frac{1-|x|^{2}}{1-\left|I_{z}(x)\right|^{2}} \cdot \frac{J_{z}(z)}{J_{z}(x)}
$$

we have $A(z)=A_{1}(z)+A_{2}(z)$, where

$$
\begin{aligned}
& A_{1}(z)=J_{z}(z) \int_{D(z, 2 r) \cap T_{z}^{\text {mod }}}\left(\partial_{p_{I_{z}(x)}} f\right)\left(I_{z}(x)\right)\left\langle z, p_{I_{z}(x)}\right\rangle \frac{1-|x|^{2}}{(1-\langle z, x\rangle)^{3}} d v_{1}(x) \quad \text { and } \\
& A_{2}(z)=\int_{D(z, 2 r) \cap T_{z}^{\bmod }}\left(\partial_{p_{I_{z}(x)}} f\right)\left(I_{z}(x)\right)\left\langle z, p_{I_{z}(x)}\right\rangle \frac{1-\left|I_{z}(x)\right|^{2}}{\left(1-\left\langle z, I_{z}(x)\right\rangle\right)^{3}} F(z, x) J_{z}(x) d v_{1}(x) .
\end{aligned}
$$

Let us first consider $A_{1}(z)$.
There is a $0<\tau_{4} \leq \tau_{3}$ such that if $0<t \leq \tau_{4}$, then for each $z \in M^{(t)}, D(z, 3 r) \subset$ $B\left(z, \min \left\{c_{1}, C_{1}^{-1} \epsilon\right\}\right)$, where $\epsilon$ is the constant in Lemma 3.6. Recall that $I_{z}(z)=z$. Thus by (3.10) and Lemma 3.6, if $0<t \leq \tau_{4}$ and $z \in M^{(t)}$, then there is a $y(z)$ in the finite set $F$ such that $I_{z}\left(D(z, 3 r) \cap T^{\text {mod }}\right) \subset U_{y(z)}$. Applying (3.7), we now have

$$
\begin{equation*}
p_{I_{z}(x)}=s_{y(z)}\left(I_{z}(x)\right) \eta_{y(z)}\left(I_{z}(x)\right) \tag{3.12}
\end{equation*}
$$

for every $x \in D(z, 3 r) \cap T_{z}^{\text {mod }}$. Define

$$
\begin{aligned}
\lambda_{z}(x) & =s_{y(z)}\left(I_{z}(x)\right)\left\langle z, p_{I_{z}(x)}\right\rangle-s_{y(z)}(z)\left\langle z, p_{z}\right\rangle \\
& =\left|s_{y(z)}\left(I_{z}(x)\right)\right|^{2}\left\langle z, \eta_{y(z)}\left(I_{z}(x)\right)\right\rangle-\left|s_{y(z)}(z)\right|^{2}\left\langle z, \eta_{y(z)}(z)\right\rangle
\end{aligned}
$$

$x \in D(z, 2 r) \cap T_{z}^{\text {mod }}$. By the uniform continuity of $\eta_{y}$ and $s_{y}$ on $\overline{U_{y}}, y \in F$, if we denote

$$
\delta(r, t)=\sup _{z \in M^{(t)}} \sup _{x \in D(z, 2 r) \cap T_{z}^{\bmod }}\left|\lambda_{z}(x)\right|,
$$

then

$$
\begin{equation*}
\lim _{t \downarrow 0} \delta(r, t)=0 \tag{3.13}
\end{equation*}
$$

for every given $1 \leq r<\infty$.
By (3.12) and the definition of $\lambda_{z}(x)$, we have $A_{1}(z)=A_{11}(z)+A_{12}(z)$, where

$$
\begin{aligned}
& A_{11}(z)=J_{z}(z) s_{y(z)}(z)\left\langle z, p_{z}\right\rangle \int_{D(z, 2 r) \cap T_{z}^{\bmod }}\left(\partial_{\eta_{y(z)}\left(I_{z}(x)\right)} f\right)\left(I_{z}(x)\right) \frac{1-|x|^{2}}{(1-\langle z, x\rangle)^{3}} d v_{1}(x) \quad \text { and } \\
& A_{12}(z)=J_{z}(z) \int_{D(z, 2 r) \cap T_{z}^{\bmod }} \lambda_{z}(x)\left(\partial_{\eta_{y(z)}\left(I_{z}(x)\right)} f\right)\left(I_{z}(x)\right) \frac{1-|x|^{2}}{(1-\langle z, x\rangle)^{3}} d v_{1}(x)
\end{aligned}
$$

Being a local inverse of $P_{z}$, the map $I_{z}$ is analytic. Thus the map $x \mapsto \eta_{y(z)}\left(I_{z}(x)\right)$ is analytic on $D(z, 3 r) \cap T_{z}^{\bmod }$. Therefore it follows from Lemma 2.10 that

$$
\begin{aligned}
A_{11}(z) & =\Phi(2 r) J_{z}(z) s_{y(z)}(z)\left\langle z, p_{z}\right\rangle\left(\partial_{\eta_{y(z)}\left(I_{z}(z)\right)} f\right)\left(I_{z}(z)\right) \\
& =\Phi(2 r) J_{z}(z) s_{y(z)}(z)\left\langle z, p_{z}\right\rangle\left(\partial_{\eta_{y(z)}(z)} f\right)(z) \\
& =\Phi(2 r) J_{z}(z)\left\langle z, p_{z}\right\rangle\left(\partial_{p_{z}} f\right)(z),
\end{aligned}
$$

where the last $=$ follows from (3.12). Recalling (2.2), (2.9) and (2.7), we see that there is a $0<C_{2}<\infty$ such that

$$
\begin{equation*}
\left|\left(\partial_{p_{z}} f\right)(z)\right| \leq C_{2}\left|\left\langle p_{z}, z\right\rangle A_{11}(z)\right| \tag{3.14}
\end{equation*}
$$

for $z \in M^{(t)}, 0<t \leq \tau_{4}$.
By Lemma 2.11, there is a $0<\tau_{5} \leq \tau_{4}$ such that if $z \in M^{\left(\tau_{5}\right)}$, then

$$
1-|x|^{2} \leq 2\left(1-\left|I_{z}(x)\right|^{2}\right) \quad \text { for every } \quad x \in D(z, 2 r) \cap T_{z}^{\bmod }
$$

Thus, applying (2.7) and the bounds for $\left|s_{y}\right|$ on $\overline{U_{y}}$, for $z \in M^{\left(\tau_{5}\right)}$ we have

$$
\begin{aligned}
\left|A_{12}(z)\right| & \leq C_{3} \delta(r, t) \int_{D(z, 2 r) \cap T_{z}^{\text {mod }}}\left|\left(\partial_{\eta_{y(z)}\left(I_{z}(x)\right)} f\right)\left(I_{z}(x)\right)\right| \frac{1-|x|^{2}}{|1-\langle z, x\rangle|^{3}} d v_{1}(x) \\
& \leq C_{4} \delta(r, t) \int_{D(z, 2 r) \cap T_{z}^{\text {mod }}}\left|\left(\partial_{\eta_{y(z)}\left(I_{z}(x)\right)} f\right)\left(I_{z}(x)\right)\right| \frac{1-\left|I_{z}(x)\right|^{2}}{\left|1-\left\langle z, I_{z}(x)\right\rangle\right|^{3}} J_{z}(x) d v_{1}(x) \\
& =C_{4} \delta(r, t) \int_{U(z)}\left|\left(\partial_{\eta_{y(z)}(w)} f\right)(w)\right| \frac{1-|w|^{2}}{|1-\langle z, w\rangle|^{3}} d v_{M}(w) \\
& \leq C_{5} \delta(r, t) \int_{M^{\left(t_{0}\right)}}\left|\left(\partial_{p_{w}} f\right)(w)\right| \frac{1}{|1-\langle z, w\rangle|^{3}} d \mu(w)
\end{aligned}
$$

Using the operator $B$ in Lemma 3.8, for $0<t \leq \tau_{5}$ we have

$$
\begin{align*}
\int_{M^{(t)}}\left|A_{12}(z)\right|^{2} d \mu(z) & \leq\left\{C_{5} \delta(r, t)\|B\|\right\}^{2} \int_{M^{\left(t_{0}\right)}}\left|\left(\partial_{p_{w}} f\right)(w)\right|^{2} d \mu(w) \\
& =\left\{C_{5} \delta(r, t)\|B\|\right\}^{2}\|f\|_{*^{2}}^{2} \tag{3.15}
\end{align*}
$$

Denote

$$
\epsilon(r, t)=\sup _{z \in M^{(t)}}\left\{\sup _{x \in D(z, 2 r) \cap T_{z}^{\bmod }}|F(z, x)|\right\}
$$

$0<t \leq \tau_{3}$. By the definition of $F(z, x)$, we can rewrite it in the form

$$
F(z, x)=\left(1-\frac{1-|x|^{2}}{1-\left|I_{z}(x)\right|^{2}}\right)+\frac{1-|x|^{2}}{1-\left|I_{z}(x)\right|^{2}} \cdot \frac{1}{J_{z}(x)} \cdot\left(J_{z}(x)-J_{z}(z)\right)
$$

From (2.7) and (2.9) we obtain an upper bound for $1 / J_{z}(x)$. Therefore it follows from Lemma 2.11 that

$$
\begin{equation*}
\lim _{t \downarrow 0} \epsilon(r, t)=0 \tag{3.16}
\end{equation*}
$$

for every $1 \leq r<\infty$. Applying (2.10) again, we have

$$
\begin{aligned}
\left|A_{2}(z)\right| & \leq \epsilon(r, t) \int_{D(z, 2 r) \cap T_{z}^{\bmod }}\left|\left(\partial_{p_{I_{z}(x)}} f\right)\left(I_{z}(x)\right)\right| \frac{1-\left|I_{z}(x)\right|^{2}}{\left|1-\left\langle z, I_{z}(x)\right\rangle\right|^{3}} J_{z}(x) d v_{1}(x) \\
& \leq \epsilon(r, t) \int_{M^{\left(t_{0}\right)}}\left|\left(\partial_{p_{w}} f\right)(w)\right| \frac{1}{|1-\langle z, w\rangle|^{3}} d \mu(w)
\end{aligned}
$$

for $z \in M^{(t)}, 0<t \leq \tau_{3}$. Thus it follows from Lemma 3.8 that if $0<t \leq \tau_{3}$, then

$$
\begin{align*}
\int_{M^{(t)}}\left|A_{2}(z)\right|^{2} d \mu(z) & \leq\{\epsilon(r, t)\|B\|\}^{2} \int_{M^{\left(t_{0}\right)}}\left|\left(\partial_{p_{w}} f\right)(w)\right|^{2} d \mu(w) \\
& =\{\epsilon(r, t)\|B\|\}^{2}\|f\|_{*}^{2} \tag{3.17}
\end{align*}
$$

For $1 \leq r<\infty$ and $0<t \leq \tau_{3}$, we define

$$
\sigma(r, t)=\sup _{z \in M^{(t)}} \sup \left\{|w-z|: w \in I_{z}\left(D(z, 2 r) \cap T_{z}^{\bmod }\right)\right\}
$$

Recalling (2.7), we have

$$
\begin{equation*}
\lim _{t \downarrow 0} \sigma(r, t)=0 \tag{3.18}
\end{equation*}
$$

for any given $1 \leq r<\infty$. By (3.11), we have

$$
|B(z)| \leq \sigma(r, t) \int_{U(z)}\left|\left(\partial_{p_{w}} f\right)(w)\right| \frac{1-|w|^{2}}{|1-\langle z, w\rangle|^{3}} d v_{M}(w) \leq \sigma(r, t) \int_{M^{\left(t_{0}\right)}} \frac{\left|\left(\partial_{p_{w}} f\right)(w)\right|}{|1-\langle z, w\rangle|^{3}} d \mu(w)
$$

for $z \in M^{(t)}, 0<t \leq \tau_{3}$. Applying Lemma 3.8, for each $0<t \leq \tau_{3}$ we now have

$$
\begin{align*}
\int_{M^{(t)}}|B(z)|^{2} d \mu(z) & \leq\{\sigma(r, t)\|B\|\}^{2} \int_{M^{\left(t_{0}\right)}}\left|\left(\partial_{p_{w}} f\right)(w)\right|^{2} d \mu(w) \\
& =\{\sigma(r, t)\|B\|\}^{2}\|f\|_{*}^{2} \tag{3.19}
\end{align*}
$$

Finally, from (3.9) we obtain

$$
|C(z)| \leq \int_{M^{\left(t_{0}\right)} \backslash D(z, r)}\left|\left(\partial_{p_{w}} f\right)(w)\right| \frac{1}{|1-\langle z, w\rangle|^{3}} d \mu(w)
$$

$z \in M^{\left(\tau_{3}\right)}$. Using the operator $B_{r}$ in Lemma 3.8, for $0<t \leq \tau_{3}$ we have

$$
\begin{equation*}
\int_{M^{(t)}}|C(z)|^{2} d \mu(z) \leq\left\|B_{r}\right\|^{2} \int_{M^{\left(t_{0}\right)}}\left|\left(\partial_{p_{w}} f\right)(w)\right|^{2} d \mu(w)=\left\|B_{r}\right\|^{2}\|f\|_{*}^{2} \tag{3.20}
\end{equation*}
$$

Retracing the above steps, we have

$$
\left\langle p_{z}, z\right\rangle A_{11}(z)=(X f)(z)-\left\langle p_{z}, z\right\rangle\left(A_{12}(z)+A_{2}(z)\right)-B(z)-C(z)
$$

Thus, for $0<t \leq \tau_{5}$, it follows from (3.14), (3.15), (3.17), (3.19) and (3.20) that

$$
\begin{align*}
& \int_{M^{(t)}}\left|\left(\partial_{p_{z}} f\right)(z)\right|^{2} d \mu(z) \leq 5 C_{2}^{2} \int_{M^{(t)}}|(X f)(z)|^{2} d \mu(z) \\
& 21) \quad+5 C_{2}^{2}\left(\left\{C_{5} \delta(r, t)\|B\|\right\}^{2}+\{\epsilon(r, t)\|B\|\}^{2}+\{\sigma(r, t)\|B\|\}^{2}+\left\|B_{r}\right\|^{2}\right)\|f\|_{*}^{2} \tag{3.21}
\end{align*}
$$

Let a $\delta>0$ be given. By Lemma 3.8, we can first pick an $r \in[1, \infty)$ such that $5 C_{2}^{2}\left\|B_{r}\right\|^{2} \leq$ $\delta / 2$. With $r$ so fixed, by (3.13), (3.16) and (3.18), we can pick a $0<\tau \leq \tau_{5}$ such that

$$
5 C_{2}^{2}\left(\left\{C_{5} \delta(r, t)\|B\|\right\}^{2}+\{\epsilon(r, t)\|B\|\}^{2}+\{\sigma(r, t)\|B\|\}^{2}\right) \leq \delta / 2
$$

for every $0<t \leq \tau$. Substituting these bounds in (3.21), the lemma is proved.
For each $f \in H_{n}^{2}$, define

$$
(Y f)(z)=\int_{M^{\left(t_{0}\right)}} \frac{\left(\partial_{p_{w}} f\right)(w)\left\langle p_{z}, p_{w}\right\rangle}{(1-\langle z, w\rangle)^{2}} d \mu(w), \quad z \in \mathbf{B} .
$$

Lemma 3.10. Given any $\delta>0$, there is a $0<\rho<t_{0}$ such that

$$
\begin{equation*}
\int_{M^{(t)}}|(Y f)(z)|^{2} d \mu(z) \leq \delta \int_{M^{\left(t_{0}\right)}}\left|\left(\partial_{p_{w}} f\right)(w)\right|^{2} d \mu(z) \tag{3.22}
\end{equation*}
$$

for all $0<t \leq \rho$ and $f \in H_{n}^{2}$.
Proof. On the Hilbert space $L^{2}(M, d \mu)$, define the operator

$$
(L \varphi)(z)=\int_{M^{\left(t_{0}\right)}} \frac{\varphi(w)\left\langle p_{z}, p_{w}\right\rangle}{(1-\langle z, w\rangle)^{2}} d \mu(w)
$$

$\varphi \in L^{2}(M, d \mu)$. It follows from Lemma 2.8 that

$$
\iint \frac{1}{|1-\langle z, w\rangle|^{4}} d \mu(w) d \mu(z)=\iint \frac{\left(1-|z|^{2}\right)\left(1-|w|^{2}\right)}{|1-\langle z, w\rangle|^{4}} d v_{M}(w) d v_{M}(z)<\infty
$$

Hence $L$ is a Hilbert-Schmidt operator on $L^{2}(M, d \mu)$. Thus for any given $\delta>0$, there is a $0<\rho<t_{0}$ such that $\left\|M_{\chi_{M(\rho)}} L\right\| \leq \delta^{1 / 2}$, where $M_{\chi_{M}(\rho)}$ is the operator of multiplication by the function $\chi_{M^{(\rho)}}$ on $L^{2}(M, d \mu)$. Obviously, (3.22) holds for this $\rho$.

Corollary 3.11. Given any $\delta>0$, there exist $0<t<t_{0}$ and $0<C<\infty$ such that

$$
\begin{equation*}
\int_{M^{(t)}}\left|\left(\partial_{p_{z}} f\right)(z)\right|^{2} d \mu(z) \leq C\left\|T_{1} f\right\|_{*}^{2}+\delta\|f\|_{*}^{2} \tag{3.23}
\end{equation*}
$$

for every $f \in H_{n}^{2}$.
Proof. Let $\delta>0$ be given. Applying Lemma 3.9 to $\delta / 2$, we obtain constants $0<\tau \leq t_{0}$ and $0<C<\infty$ such that

$$
\begin{equation*}
\int_{M^{(t)}}\left|\left(\partial_{p_{z}} f\right)(z)\right|^{2} d \mu(z) \leq C \int_{M^{(t)}}|(X f)(z)|^{2} d \mu(z)+\frac{\delta}{2}\|f\|_{*}^{2} \tag{3.24}
\end{equation*}
$$

for all $0<t \leq \tau$ and $f \in H_{n}^{2}$. For each $f \in H_{n}^{2}$, we have

$$
\left(T_{1} f\right)(z)=\int_{M^{\left(t_{0}\right)}} \frac{\left(\partial_{p_{w}} f\right)(w)\left\langle z, p_{w}\right\rangle}{(1-\langle z, w\rangle)^{2}} d \mu(w), \quad z \in \mathbf{B}
$$

By straightforward differentiation,

$$
\left(\partial_{p_{z}} T_{1} f\right)(z)=2(X f)(z)+(Y f)(z)
$$

Thus from (3.24) we obtain the inequality

$$
\begin{equation*}
\int_{M^{(t)}}\left|\left(\partial_{p_{z}} f\right)(z)\right|^{2} d \mu(z) \leq C\left\|T_{1} f\right\|_{*}^{2}+C \int_{M^{(t)}}|(Y f)(z)|^{2} d \mu(z)+\frac{\delta}{2}\|f\|_{*}^{2} \tag{3.25}
\end{equation*}
$$

for all $0<t \leq \tau$ and $f \in H_{n}^{2}$. By Lemma 3.10, there is a $0<\rho<t_{0}$ such that

$$
\begin{equation*}
C \int_{M^{(t)}}|(Y f)(z)|^{2} d \mu(z) \leq \frac{\delta}{2}\|f\|_{*}^{2} \tag{3.26}
\end{equation*}
$$

for all $0<t \leq \rho$ and $f \in H_{n}^{2}$. Combining (3.25) and (3.26), we see that (3.23) holds for every $0<t \leq \min \{\tau, \rho\}$.

Lemma 3.12. There exist a finite number of open subsets $W_{1}, \ldots, W_{m}$ of $M^{\left(t_{0}\right)}$ such that

$$
W_{1} \cup \cdots \cup W_{m} \supset M^{\left(t^{*}\right)}
$$

for some $0<t^{*}<t_{0}$ and such that the following hold true for every $1 \leq j \leq m$ :
(1) $W_{j}=G_{j}\left(\left(0, c_{j}\right) \times\left(-b_{j}, b_{j}\right)\right)$, where $0<c_{j}<t_{0}, b_{j}>0$, and
$G_{j}:\left(0, c_{j}\right) \times\left(-b_{j}, b_{j}\right) \rightarrow \mathbf{C}^{n}$ is a one-to-one $C^{\infty}$ map.
(2) There are $0<\epsilon_{j} \leq M_{j}<\infty$ such that $D G_{j}$, the derivative of $G_{j}$, satisfies the inequality $\epsilon_{j} \leq\left(D G_{j}\right)^{*}(x, y)\left(D G_{j}\right)(x, y) \leq M_{j}$ for all $(x, y) \in\left(0, c_{j}\right) \times\left(-b_{j}, b_{j}\right)$.
(3) If $w=G_{j}(x, y)$ for some $(x, y) \in\left(0, c_{j}\right) \times\left(-b_{j}, b_{j}\right)$, then $x=1-|w|^{2}$. Equivalently, for each $w \in W_{j}$, there is a unique $y_{w} \in\left(-b_{j}, b_{j}\right)$ such that $w=G_{j}\left(1-|w|^{2}, y_{w}\right)$.

Proof. Consider the function $\rho(w)=1-|w|^{2}$. Since $\mathcal{M}$ intersects $\partial \mathbf{B}$ transversely, the vector $\nabla_{\mathcal{M}} \rho$ does not vanish near $\mathcal{M} \cap \partial \mathbf{B}$. Thus we can use $\rho$ as one of the real coordinates on $\mathcal{M}$ near $\partial \mathbf{B}$. More precisely, if $\zeta \in \mathcal{M} \cap \partial \mathbf{B}$, then $\zeta$ has an open neighborhood $N_{\zeta}$ in $\mathcal{M}$ that has the following properties:
$(\alpha) N_{\zeta}=G((-c, c) \times(-b, b))$, where $0<c<t_{0}, b>0$ and
$G:(-c, c) \times(-b, b) \rightarrow \mathbf{C}^{n}$ is a one-to-one $C^{\infty}$ map.
( $\beta$ ) There are $0<\epsilon \leq M<\infty$ such that $D G$, the derivative of $G$, satisfies the matrix inequality $\epsilon \leq(D G)^{*}(x, y)(D G)(x, y) \leq M$ for all $(x, y) \in(-c, c) \times(-b, b)$.
$(\gamma)$ If $w=G(x, y)$ for some $(x, y) \in(-c, c) \times(-b, b)$, then $x=1-|w|^{2}$. Equivalently, for each $w \in N_{\zeta}$, there is a unique $y_{w} \in(-b, b)$ such that $w=G\left(1-|w|^{2}, y_{w}\right)$.
We then define $W_{\zeta}=N_{\zeta} \cap \mathbf{B}$. By $(\gamma)$ and $(\alpha)$, we have $W_{\zeta}=G((0, c) \times(-b, b)) \subset M^{\left(t_{0}\right)}$. Since $\mathcal{M} \cap \partial \mathbf{B}$ is compact, there is a finite subset $Z$ of $\mathcal{M} \cap \partial \mathbf{B}$ such that $\cup_{\zeta \in Z} N_{\zeta} \supset \mathcal{M} \cap \partial \mathbf{B}$. Since $\cup_{\zeta \in Z} N_{\zeta}$ is an open subset of $\mathcal{M}$, there is a $0<t^{*}<t_{0}$ such that $\cup_{\zeta \in Z} N_{\zeta} \supset M^{\left(t^{*}\right)}$. Obviously, if we re-enumerate the finite family of sets $\left\{W_{\zeta}: \zeta \in Z\right\}$ as $\left\{W_{1}, \ldots, W_{m}\right\}$, then the lemma holds.

Proposition 3.13. The dimension of $\left\{f \in \mathcal{Q}: T_{1} f=0\right\}$ is finite.
Proof. Let $g \in H_{n}^{2}$. If $T_{1} g=0$, then

$$
\int_{M^{\left(t_{0}\right)}}\left|\left(\partial_{p_{w}} g\right)(w)\right|^{2} d \mu(w)=0
$$

Let $y \in M^{\left(t_{0}\right)}$. By the argument in the proof of Lemma 3.7 (see (3.8)), the above equality implies that $\left(g \circ \rho_{y}\right)^{\prime}=0$ on $B_{1}(0, r)$ for some $r=r(y)>0$. That is, $g \circ \rho_{y}$ is a constant on $B_{1}(0, r)$. Equivalently, $g$ is a constant on an open subset of $M^{\left(t_{0}\right)}$ containing $y$.

Let $W_{1}, \ldots, W_{m}$ be the open subsets of $M^{\left(t_{0}\right)}$ provided by Lemma 3.12. For each $1 \leq j \leq m$, since $W_{j}$ is homeomorphic to $\left(0, c_{j}\right) \times\left(-b_{j}, b_{j}\right)$, it is a connected subset of $M^{\left(t_{0}\right)}$. Thus by the conclusion of the preceding paragraph, if $g \in H_{n}^{2}$ and $T_{1} g=0$, then $g$ is a constant on $W_{j}$ for every $1 \leq j \leq m$.

Thus we can define a linear map $L:\left\{f \in \mathcal{Q}: T_{1} f=0\right\} \rightarrow \mathbf{C}^{m}$ by the formula

$$
L f=\left(f\left|W_{1}, \ldots, f\right| W_{m}\right)
$$

where $f \mid W_{j}$ means the constant value of $f$ on $W_{j}, 1 \leq j \leq m$. If $h$ is in the kernel of $L$, then $h=0$ on $W_{1} \cup \cdots \cup W_{m}$. By the fact $W_{1} \cup \cdots \cup W_{m} \supset M^{\left(t^{*}\right)}$ and the maximum modulus principle [8, pages 72,73], we have $h=0$ on $M$. That is, $h \perp \mathcal{Q}$. Since $h \in \mathcal{Q}$, this means $h=0$. Hence $\operatorname{dim}\left\{f \in \mathcal{Q}: T_{1} f=0\right\} \leq m$.

Proof of Theorem 3.5. Let $t^{*} \in\left(0, t_{0}\right)$ be the number provided by Lemma 3.12. By Corollary 3.11, there are $0<t<t^{*}$ and $0<C<\infty$ such that

$$
\int_{M^{(t)}}\left|\left(\partial_{p_{z}} f\right)(z)\right|^{2} d \mu(z) \leq C\left\|T_{1} f\right\|_{*}^{2}+\frac{1}{2}\|f\|_{*}^{2}
$$

for every $f \in H_{n}^{2}$. After the obvious cancellation, we obtain the inequality

$$
\begin{equation*}
\frac{1}{2} \int_{M^{(t)}}\left|\left(\partial_{p_{z}} f\right)(z)\right|^{2} d \mu(z) \leq C\left\|T_{1} f\right\|_{*}^{2}+\frac{1}{2} \int_{M^{\left(t_{0}\right)} \backslash M^{(t)}}\left|\left(\partial_{p_{z}} f\right)(z)\right|^{2} d \mu(z) \tag{3.27}
\end{equation*}
$$

for every $f \in H_{n}^{2}$.
Pick a positive number $t_{1}>0$ satisfying the condition $t_{1}<\min \left\{t, c_{1}, \ldots, c_{m}\right\}$, where $c_{1}, \ldots, c_{m}$ are the same as in Lemma 3.12. With this $t_{1}$, we define the operator

$$
\Phi=\int_{N^{\left(t_{1}\right)}} K_{w} \otimes K_{w} d \mu(w)
$$

(see (2.13)). Then $\Phi$ is obviously a positive operator, and we have

$$
\operatorname{tr}(\Phi)=\int_{N^{\left(t_{1}\right)}} \frac{1}{1-|w|^{2}} d \mu(w) \leq \int_{M} 1 d v_{M}(w)<\infty
$$

Thus $\Phi$ is in the trace class, but here we only need the compactness of $\Phi$. Define

$$
S_{1}=T_{1}+\Phi
$$

We claim that there is an $a>0$ such that the spectrum of $S_{1}$ does not intersect the interval $(0, a)$. Postponing the proof of this claim for a moment, we first show that this claim implies the conclusion of Theorem 3.5.

Of course, both $T_{1}$ and $\Phi$ map $\mathcal{Q}$ into itself. Since $T_{1}$ and $\Phi$ are both positive, we have $\left\{f \in \mathcal{Q}: S_{1} f=0\right\} \subset\left\{f \in \mathcal{Q}: T_{1} f=0\right\}$. Thus it follows from Proposition 3.13 that $\operatorname{dim}\left\{f \in \mathcal{Q}: S_{1} f=0\right\}<\infty$. If the spectrum of $S_{1}$ does not intersect $(0, a)$ for some $a>0$, then 0 is not in the essential spectrum of the restricted operator $S_{1} \mid \mathcal{Q}$. Since $\Phi$ is compact, this means that 0 is not in the essential spectrum of the restricted operator $T_{1} \mid \mathcal{Q}$. Thus there is a $c>0$ such that the spectrum of $T_{1} \mid \mathcal{Q}$ does not intersect $(0, c)$. Since $T_{1}=0$ on $\mathcal{Q}^{\perp}$, it follows that the spectrum of $T_{1}$ does not intersect $(0, c)$.

Thus we have reduced the proof of Theorem 3.5 to the proof of the claim that there is an $a>0$ such that the spectrum of $S_{1}$ does not intersect $(0, a)$. To prove this claim, let $d E$ be the spectral measure for the positive operator $S_{1}$. That is,

$$
S_{1}=\int_{0}^{\left\|S_{1}\right\|} \lambda d E(\lambda)
$$

Suppose that $E(0, a) \neq 0$ for every $a>0$. We will complete the proof by showing that this leads to a contradiction. For each $k \in \mathbf{N}$, since $E(0,1 / k) \neq 0$, we pick an $f_{k} \in E(0,1 / k) H_{n}^{2}$ such that $\left\langle S_{1} f_{k}, f_{k}\right\rangle=1$. That is,

$$
\begin{equation*}
\int_{M^{\left(t_{0}\right)}}\left|\left(\partial_{p_{w}} f_{k}\right)(w)\right|^{2} d \mu(w)+\int_{N^{\left(t_{1}\right)}}\left|f_{k}(w)\right|^{2} d \mu(w)=1 \tag{3.28}
\end{equation*}
$$

for every $k$. Obviously, the sequence $\left\{S_{1}^{1 / 2} f_{k}\right\}$ weakly converges to 0 in $H_{n}^{2}$. Since $T_{1} \leq S_{1}$, there is a contraction $A$ such that $T_{1}^{1 / 2}=A S_{1}^{1 / 2}$. Hence the sequence $\left\{T_{1}^{1 / 2} f_{k}\right\}$ also weakly converges to 0 .

Let $0<\epsilon<\min \left\{t_{1}, t_{0}-t_{1}\right\}$. Then the closure of $N^{\left(t_{1}-\epsilon\right)} \cap M^{\left(t_{1}+\epsilon\right)}$ is a compact subset of $M^{\left(t_{0}\right)}$. By Lemma 3.7 and a usual covering argument, the weak convergence to 0 of the sequence $\left\{T_{1}^{1 / 2} f_{k}\right\}$ implies

$$
\begin{equation*}
\lim _{k \rightarrow \infty} \sup \left\{\left|\left(\partial_{p_{z}} f_{k}\right)(z)\right|: z \in N^{\left(t_{1}-\epsilon\right)} \cap M^{\left(t_{1}+\epsilon\right)}\right\}=0 \tag{3.29}
\end{equation*}
$$

Denote $\Delta=\left\{z \in M: 1-|z|^{2}=t_{1}+(\epsilon / 2)\right\}$. By the choices of $t_{1}$ and $\epsilon$ above, $\Delta$ is a compact set in $\left\{w \in M: 1-|w|^{2}>t_{1}\right\}$. Thus $\Delta$ can be covered by a finite number of open sets $D_{1}, \ldots, D_{\ell}$ in $\left\{w \in M: 1-|w|^{2}>t_{1}\right\}$ in such a way that each $D_{i}$ is biholomorphically equivalent to the unit disc $D=\{\xi \in \mathbf{C}:|\xi|<1\}$. By the Bergman integral formula, there is a constant $0<C_{1}<\infty$ such that

$$
\begin{equation*}
\sup _{z \in \Delta}|g(z)|^{2} \leq C_{1} \int_{N^{\left(t_{1}\right)}}|g(w)|^{2} d \mu(w) \tag{3.30}
\end{equation*}
$$

for every $g \in H_{n}^{2}$. Combining this with (3.28), we see that

$$
\begin{equation*}
\sup _{z \in \Delta}\left|f_{k}(z)\right|^{2} \leq C_{1} \quad \text { for every } \quad k \tag{3.31}
\end{equation*}
$$

For $k \in \mathbf{N}, 1 \leq j \leq m, s \in\left(-b_{j}, b_{j}\right)$ and $u \in\left[t_{1}-(\epsilon / 2), t_{1}+(\epsilon / 2)\right]$, we can write

$$
\begin{equation*}
f_{k}\left(G_{j}(u, s)\right)=f_{k}\left(G_{j}\left(t_{1}+(\epsilon / 2), s\right)\right)-\int_{u}^{t_{1}+(\epsilon / 2)} \frac{d}{d r} f_{k}\left(G_{j}(r, s)\right) d r \tag{3.32}
\end{equation*}
$$

where $G_{j}$ and $b_{j}$ are the same as in Lemma 3.12. Taking the derivative in the integral, combining (3.32) with (3.31), (3.29) and Lemma 3.12, and using the fact $\operatorname{dim}_{\mathbf{C}} T_{z}=1$ and lower bound (2.2), we deduce that there is a $0<C_{2}<\infty$ such that

$$
\sup \left\{\left|f_{k}(z)\right|: z \in N^{\left(t_{1}-(\epsilon / 2)\right)} \cap M^{\left(t_{1}+(\epsilon / 2)\right)}\right\} \leq C_{2} \quad \text { for every } k
$$

By the maximum modulus principle, this implies that

$$
\begin{equation*}
\sup \left\{\left|f_{k}(z)\right|: z \in N^{\left(t_{1}-(\epsilon / 2)\right)}\right\} \leq C_{2} \quad \text { for every } k \tag{3.33}
\end{equation*}
$$

If $\varphi$ is a bounded analytic function on $D$, then

$$
\varphi^{\prime}(0)=\frac{1}{\pi} \int_{-\pi}^{\pi} \int_{0}^{1} \varphi\left(r e^{i \theta}\right) e^{-i \theta} d r d \theta
$$

Using this identity and (3.33), and using the bounds for $s_{y}, y \in F$, again, we obtain a $0<C_{3}<\infty$ such that

$$
\begin{equation*}
\sup \left\{\left|\left(\partial_{p_{z}} f_{k}\right)(z)\right|: z \in N^{\left(t_{1}-(\epsilon / 3)\right)} \cap M^{\left(t_{0}\right)}\right\} \leq C_{3} \quad \text { for every } k \tag{3.34}
\end{equation*}
$$

Since the sequence $\left\{T_{1}^{1 / 2} f_{k}\right\}$ weakly converges to 0 , Lemma 3.7 tells us that

$$
\lim _{k \rightarrow \infty}\left(\partial_{p_{w}} f_{k}\right)(w)=0 \quad \text { for every } \quad w \in M^{\left(t_{0}\right)}
$$

Combining this pointwise convergence with (3.34) and with the dominated convergence theorem, we have

$$
\begin{equation*}
\lim _{k \rightarrow \infty} \int_{M^{\left(t_{0}\right)} \backslash M^{(t)}}\left|\left(\partial_{p_{w}} f_{k}\right)(w)\right|^{2} d \mu(w)=0 . \tag{3.35}
\end{equation*}
$$

Similar to (3.30), for each $t_{1}<u<t_{0}$ there is a $0<C(u)<\infty$ such that

$$
\sup \left\{|g(z)|^{2}: z \in M \text { and } 1-|z|^{2}=u\right\} \leq C(u) \int_{N^{\left(t_{1}\right)}}|g(w)|^{2} d \mu(w) \quad \text { for every } g \in H_{n}^{2}
$$

Again by the maximum modulus principle, the above implies that for each $t_{1}<u<t_{0}$,

$$
\sup \left\{|g(z)|^{2}: z \in M \text { and } 1-|z|^{2} \geq u\right\} \leq C(u) \int_{N^{\left(t_{1}\right)}}|g(w)|^{2} d \mu(w) \quad \text { for every } g \in H_{n}^{2}
$$

Let $\mathcal{H}$ be the closure of $H_{n}^{2}$ in $L^{2}\left(N^{\left(t_{1}\right)}, d \mu\right)$. The above bound means that for each $z \in N^{\left(t_{1}\right)}$, the map $g \mapsto g(z)$ extends to a bounded linear functional on $\mathcal{H}$. Since $\left\{S_{1}^{1 / 2} f_{k}\right\}$ weakly converges to 0 and $\Phi \leq S_{1}$, the sequence $\left\{\Phi^{1 / 2} f_{k}\right\}$ also weakly converges to 0 in $H_{n}^{2}$. For any $h \in H_{n}^{2}$, we have

$$
\int_{N^{\left(t_{1}\right)}} f_{k}(w) \overline{h(w)} d \mu(w)=\left\langle\Phi^{1 / 2} f_{k}, \Phi^{1 / 2} h\right\rangle, \quad k \in \mathbf{N}
$$

Thus in the Hilbert space $\mathcal{H}$, the sequence $\left\{f_{k}\right\}$ weakly converges to 0 . Hence

$$
\lim _{k \rightarrow \infty} f_{k}(z)=0
$$

for every $z \in N^{\left(t_{1}\right)}$. Combining this pointwise convergence with (3.33) and with the dominated convergence theorem, we obtain

$$
\lim _{k \rightarrow \infty} \int_{N^{\left(t_{1}\right)}}\left|f_{k}(w)\right|^{2} d \mu(w)=0
$$

From this limit and (3.28), (3.35) it follows that

$$
\begin{equation*}
\lim _{k \rightarrow \infty} \int_{M^{(t)}}\left|\left(\partial_{p_{w}} f_{k}\right)(w)\right|^{2} d \mu(w)=1 \tag{3.36}
\end{equation*}
$$

Next we show that

$$
\begin{equation*}
\lim _{k \rightarrow \infty}\left\|T_{1} f_{k}\right\|_{*}^{2}=0 \tag{3.37}
\end{equation*}
$$

Indeed for each $k \in \mathbf{N}$, we have

$$
\begin{equation*}
\left\|T_{1} f_{k}\right\|_{*}^{2}=\left\langle T_{1} T_{1} f_{k}, T_{1} f_{k}\right\rangle=\left\langle T_{1} S_{1} f_{k}, S_{1} f_{k}\right\rangle-\left\langle T_{1} \Phi f_{k}, T_{1} f_{k}\right\rangle-\left\langle T_{1} S_{1} f_{k}, \Phi f_{k}\right\rangle \tag{3.38}
\end{equation*}
$$

Since $T_{1} \leq S_{1}, f_{k} \in E(0,1 / k) H_{n}^{2}$ and $\left\|S_{1}^{1 / 2} f_{k}\right\|=1$, we have

$$
\begin{equation*}
\left\langle T_{1} S_{1} f_{k}, S_{1} f_{k}\right\rangle \leq\left\langle S_{1} S_{1} f_{k}, S_{1} f_{k}\right\rangle=\left\|S_{1}^{3 / 2} f_{k}\right\|^{2} \leq k^{-2}\left\|S_{1}^{1 / 2} f_{k}\right\|^{2}=k^{-2} \tag{3.39}
\end{equation*}
$$

Since the sequence $\left\{\Phi^{1 / 2} f_{k}\right\}$ weakly converges to 0 and $\Phi^{1 / 2}$ is compact, the sequence $\left\{\Phi f_{k}\right\}$ converges to 0 strongly, i.e.,

$$
\begin{equation*}
\lim _{k \rightarrow \infty}\left\|\Phi f_{k}\right\|=0 \tag{3.40}
\end{equation*}
$$

By (3.39), the first term on the right-hand side of (3.38) tends to 0 as $k \rightarrow \infty$. Since $\left\|T_{1} f_{k}\right\|$ $\leq\left\|T_{1}^{1 / 2}\right\|\left\|T_{1}^{1 / 2} f_{k}\right\| \leq\left\|T_{1}^{1 / 2}\right\|$, it follows from (3.40) that the second term on the right-hand side of (3.38) tends to 0 as $k \rightarrow \infty$. Since $\left\|S_{1} f_{k}\right\| \leq\left\|S_{1}^{1 / 2}\right\|\left\|S_{1}^{1 / 2} f_{k}\right\|=\left\|S_{1}^{1 / 2}\right\|$, (3.40) also implies that the third term on the right-hand side of (3.38) tends to 0 as $k \rightarrow \infty$. This proves (3.37).

Now, recalling (3.27), we have

$$
\begin{equation*}
\frac{1}{2} \int_{M^{(t)}}\left|\left(\partial_{p_{z}} f_{k}\right)(z)\right|^{2} d \mu(z) \leq C\left\|T_{1} f_{k}\right\|_{*}^{2}+\frac{1}{2} \int_{M^{\left(t_{0}\right)} \backslash M^{(t)}}\left|\left(\partial_{p_{z}} f_{k}\right)(z)\right|^{2} d \mu(z) \tag{3.41}
\end{equation*}
$$

for every $k \in \mathbf{N}$. But the combination of (3.41), (3.35), (3.36) and (3.37) gives us the contradiction $1 / 2 \leq 0$. This proves our claim that there is an $a>0$ such that the spectrum of $S_{1}$ does not intersect $(0, a)$, which in turn completes the proof of Theorem 3.5.

## 4. The range space

In addition to the operator $T_{1}$ given by (3.6), we define the operator

$$
\begin{equation*}
T_{2}=\int_{M^{\left(t_{0}\right)}} K_{w} \otimes K_{w} d \mu(w) \tag{4.1}
\end{equation*}
$$

Again, $T_{2}$ is a positive operator on the Drury-Arveson space $H_{n}^{2}$. We have

$$
\operatorname{tr}\left(T_{2}\right)=\int_{M^{\left(t_{0}\right)}} \frac{1}{1-|w|^{2}} d \mu(w)=\int_{M^{\left(t_{0}\right)}} 1 d v_{M}(w)<\infty
$$

This shows that $T_{2}$ belongs to the trace class $\mathcal{C}_{1}$. Obviously, $T_{2}$ maps $H_{n}^{2}$ into $\mathcal{Q}$. We now define the operator

$$
\begin{equation*}
T=T_{1}+T_{2} \tag{4.2}
\end{equation*}
$$

on $H_{n}^{2}$. Then $T$ maps $H_{n}^{2}$ into $\mathcal{Q}$.

Proposition 4.1. There is a $c^{\prime}>0$ such that the spectrum of $T$ does not intersect the interval ( $0, c^{\prime}$ ). Moreover, $Q$ equals the spectral projection of $T$ corresponding to the interval $\left[c^{\prime}, \infty\right)$.
Proof. It will be convenient to use the following notation for this proof: If $A$ is an operator on $H_{n}^{2}$ such that $A \mathcal{Q} \subset \mathcal{Q}$, we write $A \mid \mathcal{Q}$ for the restriction of $A$ to $\mathcal{Q}$.

Theorem 3.5 tells us that the spectrum of $T_{1}$ does not intersect $(0, c)$ for some $c>0$. By Proposition 3.13, $\operatorname{dim}\left\{f \in \mathcal{Q}: T_{1} f=0\right\}<\infty$. Thus 0 is not in the essential spectrum of $T_{1} \mid \mathcal{Q}$. Since $T_{2} \mid \mathcal{Q}$ is in the trace class, we conclude that 0 is not in the essential spectrum of $T \mid \mathcal{Q}$. Consequently, there is a $c^{\prime}>0$ such that the spectrum of $T \mid \mathcal{Q}$ does not intersect the interval $\left(0, c^{\prime}\right)$. With respect to the orthogonal decomposition $H_{n}^{2}=\mathcal{R} \oplus \mathcal{Q}$, we have $T=0 \oplus(T \mid \mathcal{Q})$. Therefore the spectrum of $T$ does not intersect $\left(0, c^{\prime}\right)$. Once this fact is established, the assertion that $Q$ equals the spectral projection of $T$ corresponding to the interval $\left[c^{\prime}, \infty\right)$ is equivalent to the assertion that $\operatorname{ker}(T \mid \mathcal{Q})=\{0\}$.

To prove this last assertion, let $f \in \mathcal{Q}$ be such that $T f=0$. Since both $T_{1}$ and $T_{2}$ are positive operators, the condition $T f=0$ implies $T_{2} f=0$, i.e.,

$$
\int_{M^{\left(t_{0}\right)}}|f(w)|^{2} d \mu(w)=0
$$

This means that $f=0$ on $M^{\left(t_{0}\right)}$. By the maximum modulus principle, we have $f=0$ on $M$. Thus $f \perp \mathcal{Q}$. Since $f \in \mathcal{Q}, f$ is the zero element. This completes the proof.

Let $f$ be a $C^{1}$ function on an open set containing the closure of $M$. We define
$\|f\|_{\#}=\left\{\int_{M^{\left(t_{0}\right)}}\left|\left(\partial_{p_{w}} f\right)(w)\right|^{2} d \mu(w)+\int_{M^{\left(t_{0}\right)}}\left|\left(\overline{\partial_{p_{w}}} f\right)(w)\right|^{2} d \mu(w)+\int_{M^{\left(t_{0}\right)}}|f(w)|^{2} d \mu(w)\right\}^{1 / 2}$
(see (2.21)). Let $\mathcal{L}_{0}$ be the collection of all such $f$ with the property $\|f\|_{\#}<\infty$. Then $\|\cdot\|_{\#}$ is a norm on $\mathcal{L}_{0}$. This norm is designed to have the symmetric property $\|\bar{f}\|_{\#}=\|f\|_{\#}$, which will be important later on.

Obviously, the norm $\|\cdot\|_{\#}$ is induced by the inner product

$$
\begin{aligned}
\langle f, g\rangle_{\#} & =\int_{M^{\left(t_{0}\right)}}\left(\partial_{p_{w}} f\right)(w) \overline{\left(\partial_{p_{w}} g\right)(w)} d \mu(w)+\int_{M^{\left(t_{0}\right)}}\left(\overline{\partial_{p_{w}}} f\right)(w) \overline{\left(\overline{\partial_{p_{w}}} g\right)(w)} d \mu(w) \\
& +\int_{M^{\left(t_{0}\right)}} f(w) \overline{g(w)} d \mu(w),
\end{aligned}
$$

$f, g \in \mathcal{L}_{0}$. Let $\mathcal{L}$ denote the completion of $\mathcal{L}_{0}$ with respect to the norm $\|\cdot\|_{\#}$. Then $\mathcal{L}$ is a Hilbert space.

Definition 4.2. (a) Let $\mathcal{P}$ be the closure of the analytic polynomials $\mathbf{C}\left[z_{1}, \ldots, z_{n}\right]$ in $\mathcal{L}$.
(b) Let $P$ denote the orthogonal projection from $\mathcal{L}$ onto $\mathcal{P}$.

Obviously, if $f \in H_{n}^{2}$, then

$$
\begin{equation*}
\|f\|_{\#}^{2}=\langle T f, f\rangle=\left\|T^{1 / 2} f\right\|^{2} \tag{4.3}
\end{equation*}
$$

Since $\mathbf{C}\left[z_{1}, \ldots, z_{n}\right]$ is dense in $H_{n}^{2}$, every $f \in H_{n}^{2}$ is naturally an element in $\mathcal{P}$.
Definition 4.3. Let $J$ denote the operator that takes each $f \in H_{n}^{2}$ to the same $f$ in $\mathcal{P}$.
Thus we can rewrite (4.3) in the form of the operator identity

$$
\begin{equation*}
J^{*} J=T . \tag{4.4}
\end{equation*}
$$

Intuitively, we think of $J$ as restricting each $f \in H_{n}^{2}$ to the set $M^{\left(t_{0}\right)}$. We call $\mathcal{P}$ the range space for the restriction operator $J$. If $f \in \mathcal{R}$, then we obviously have $J f=0$. On the other hand, by Proposition 4.1,

$$
\|J f\|_{\#}=\left\|T^{1 / 2} f\right\| \geq \sqrt{c^{\prime}}\|f\| \quad \text { for every } \quad f \in \mathcal{Q} .
$$

Thus $J$ is an invertible operator that maps $\mathcal{Q}$ onto $\mathcal{P}$.
We define the operators

$$
\begin{aligned}
& \left(\hat{T}_{1} f\right)(z)=\int_{M^{\left(t_{0}\right)}} \frac{\left\langle z, p_{w}\right\rangle}{(1-\langle z, w\rangle)^{2}}\left(\partial_{p_{w}} f\right)(w) d \mu(w) \quad \text { and } \\
& \left(\hat{T}_{2} f\right)(z)=\int_{M^{\left(t_{0}\right)}} \frac{1}{1-\langle z, w\rangle} f(w) d \mu(w)
\end{aligned}
$$

$f \in \mathcal{L}_{0}$.
Lemma 4.4. The operators $\hat{T}_{1}$ and $\hat{T}_{2}$ are bounded on $\mathcal{L}_{0}$. Therefore $\hat{T}_{1}$ and $\hat{T}_{2}$ naturally extend to bounded operators on $\mathcal{L}$.

Proof. By Lemma 2.8 and the Schur test, the kernels

$$
\frac{\left\langle p_{z}, w\right\rangle\left\langle z, p_{w}\right\rangle}{(1-\langle z, w\rangle)^{3}}, \quad \frac{\left\langle z, p_{w}\right\rangle}{(1-\langle z, w\rangle)^{2}}, \quad \frac{\left\langle p_{z}, w\right\rangle}{(1-\langle z, w\rangle)^{2}} \quad \text { and } \frac{1}{1-\langle z, w\rangle}
$$

define bounded operators on $L^{2}\left(M^{\left(t_{0}\right)}, d \mu\right)$. Let us verify the details for the case

$$
K(z, w)=\frac{\left\langle p_{z}, w\right\rangle\left\langle z, p_{w}\right\rangle}{(1-\langle z, w\rangle)^{3}}
$$

the other cases are similar.
Define the function $h(z)=\left(1-|z|^{2}\right)^{-1 / 2}$ on $M^{\left(t_{0}\right)}$. By (3.1) and Lemma 2.8,

$$
\int_{M^{\left(t_{0}\right)}}|K(z, w)| h(w) d \mu(w) \leq \int_{M^{\left(t_{0}\right)}} \frac{\left(1-|w|^{2}\right)^{1 / 2}}{|1-\langle z, w\rangle|^{1+1+(1 / 2)+(1 / 2)}} d v_{M}(w) \leq C_{1} h(z) .
$$

Similarly,

$$
\int_{M^{\left(t_{0}\right)}}|K(z, w)| h(z) d \mu(z) \leq C_{1} h(w) .
$$

Thus by the Schur test, the kernel $K(z, w)$ represents a bounded operator on $L^{2}\left(M^{\left(t_{0}\right)}, d \mu\right)$.
Let $f \in \mathcal{L}_{0}$. Then

$$
\left(\partial_{p_{z}} \hat{T}_{1} f\right)(z)=\int_{M^{\left(t_{0}\right)}} \frac{2\left\langle p_{z}, w\right\rangle\left\langle z, p_{w}\right\rangle\left(\partial_{p_{w}} f\right)(w)}{(1-\langle z, w\rangle)^{3}} d \mu(w)+\int_{M^{\left(t_{0}\right)}} \frac{\left\langle p_{z}, p_{w}\right\rangle\left(\partial_{p_{w}} f\right)(w)}{(1-\langle z, w\rangle)^{2}} d \mu(w)
$$

Note that $\overline{\partial_{p_{z}}} \hat{T}_{1} f=0$. By the boundedness of the kernels mentioned in the first paragraph, $\hat{T}_{1}$ is bounded on $\mathcal{L}_{0}$. Similarly, $\hat{T}_{2}$ is also bounded on $\mathcal{L}_{0}$.

On the Hilbert space $\mathcal{L}$ we now define the operator

$$
\hat{T}=\hat{T}_{1}+\hat{T}_{2}
$$

Our next lemma is crucial to the proof of the 1-essential normality of $\mathcal{Q}$. It deals with a rarity in operator theory: a situation where self-adjointness is not so obvious.
Lemma 4.5. With respect to the inner product $\langle\cdot, \cdot\rangle_{\#}$, the operator $\hat{T}$ is self-adjoint.
Proof. For any $f \in \mathcal{L}_{0}$, straightforward differentiation gives us

$$
\begin{aligned}
& \left(\partial_{p_{z}} \hat{T}_{1} f\right)(z)=\int_{M^{\left(t_{0}\right)}} \frac{2\left\langle p_{z}, w\right\rangle\left\langle z, p_{w}\right\rangle\left(\partial_{p_{w}} f\right)(w)}{(1-\langle z, w\rangle)^{3}} d \mu(w)+\int_{M^{\left(t_{0}\right)}} \frac{\left\langle p_{z}, p_{w}\right\rangle\left(\partial_{p_{w}} f\right)(w)}{(1-\langle z, w\rangle)^{2}} d \mu(w), \\
& \left(\partial_{p_{z}} \hat{T}_{2} f\right)(z)=\int_{M^{\left(t_{0}\right)}} \frac{\left\langle p_{z}, w\right\rangle}{(1-\langle z, w\rangle)^{2}} f(w) d \mu(w)
\end{aligned}
$$

Also, $\left(\overline{\partial_{p_{z}}} \hat{T}_{1} f\right)(z)=0$ and $\left(\overline{\partial_{p_{z}}} \hat{T}_{2} f\right)(z)=0$. Thus for $f, g \in \mathcal{L}_{0}$, we have

$$
\begin{align*}
\langle\hat{T} f, g\rangle_{\#} & =\int_{M^{\left(t_{0}\right)}}\left(\partial_{p_{z}} \hat{T}_{1} f\right)(z) \overline{\left(\partial_{p_{z}} g\right)(z)} d \mu(z)+\int_{M^{\left(t_{0}\right)}}\left(\partial_{p_{z}} \hat{T}_{2} f\right)(z) \overline{\left(\partial_{p_{z}} g\right)(z)} d \mu(z) \\
& +\int_{M^{\left(t_{0}\right)}}\left(\hat{T}_{1} f\right)(z) \overline{g(z)} d \mu(z)+\int_{M^{\left(t_{0}\right)}}\left(\hat{T}_{2} f\right)(z) \overline{g(z)} d \mu(z) \tag{4.5}
\end{align*}
$$

By the integral formula for $\partial_{p_{z}} \hat{T}_{1} f$ and a change of order of integration, we have

$$
\begin{equation*}
\int_{M^{\left(t_{0}\right)}}\left(\partial_{p_{z}} \hat{T}_{1} f\right)(z) \overline{\left(\partial_{p_{z}} g\right)(z)} d \mu(z)=\int_{M^{\left(t_{0}\right)}}\left(\partial_{p_{w}} f\right)(w) \overline{\left(\partial_{p_{w}} \hat{T}_{1} g\right)(w)} d \mu(w) \tag{4.6}
\end{equation*}
$$

It is obvious that

$$
\begin{equation*}
\int_{M^{\left(t_{0}\right)}}\left(\hat{T}_{2} f\right)(z) \overline{g(z)} d \mu(z)=\int_{M^{\left(t_{0}\right)}} f(w) \overline{\left(\hat{T}_{2} g\right)(w)} d \mu(w) \tag{4.7}
\end{equation*}
$$

Finally, by the formula for $\partial_{p_{z}} \hat{T}_{2} f$, we have

$$
\begin{align*}
\int_{M^{\left(t_{0}\right)}}\left(\partial_{p_{z}} \hat{T}_{2} f\right)(z) \overline{\left(\partial_{p_{z}} g\right)(z)} d \mu(z) & =\int_{M^{\left(t_{0}\right)}} f(w) \overline{\left(\hat{T}_{1} g\right)(w)} d \mu(w) \quad \text { and }  \tag{4.8}\\
\int_{M^{\left(t_{0}\right)}}\left(\hat{T}_{1} f\right)(z) \overline{g(z)} d \mu(z) & =\int_{M^{\left(t_{0}\right)}}\left(\partial_{p_{w}} f\right)(w) \overline{\left(\partial_{p_{w}} \hat{T}_{2} g\right)(w)} d \mu(w) \tag{4.9}
\end{align*}
$$

Combining (4.5)-(4.9) with the fact that $\left(\overline{\partial_{p_{w}}} \hat{T} g\right)(w)=0$, we see that $\langle\hat{T} f, g\rangle_{\#}=\langle f, \hat{T} g\rangle_{\#}$. This completes the proof.

From the above proof we see that the individual operators $\hat{T}_{1}$ and $\hat{T}_{2}$ are not selfadjoint on $\mathcal{L}$. But miraculously, somehow the sum $\hat{T}=\hat{T}_{1}+\hat{T}_{2}$ is self-adjoint.

Proposition 4.6. (a) $\hat{T}$ maps $\mathcal{L}$ into $\mathcal{P}$.
(b) Let $\tilde{T}$ denote the restriction of $\hat{T}$ to the subspace $\mathcal{P}$. Then $\tilde{T}=J J^{*}$. In particular, $\tilde{T}$ is invertible on $\mathcal{P}$.
(c) With respect to the orthogonal decomposition $\mathcal{L}=\mathcal{P} \oplus \mathcal{P}^{\perp}$, we have $\hat{T}=\tilde{T} \oplus 0$.

Proof. (a) Recall that the kernels listed in the proof of Lemma 4.4 define bounded operators on $L^{2}\left(M^{\left(t_{0}\right)}, d \mu\right)$. Therefore for any $f \in \mathcal{L}_{0}$, we have

$$
\lim _{t \downarrow 0}\left\|\int_{M^{(t)}}\left(\partial_{p_{w}} f\right)(w) J K_{w, p_{w}} d \mu(w)\right\|_{\#}=0 .
$$

Since we already know that $J H_{n}^{2} \subset \mathcal{P}$, we have

$$
\int_{M^{\left(t_{0}\right)} \backslash M^{(t)}}\left(\partial_{p_{w}} f\right)(w) J K_{w, p_{w}} d \mu(w) \in \mathcal{P}
$$

for every $0<t<t_{0}$. Therefore

$$
\hat{T}_{1} f=\int_{M^{\left(t_{0}\right)}}\left(\partial_{p_{w}} f\right)(w) J K_{w, p_{w}} d \mu(w) \in \mathcal{P}
$$

A similar argument shows that $\hat{T}_{2} f \in \mathcal{P}$ for $f \in \mathcal{L}_{0}$. Thus $\hat{T} \mathcal{L}_{0} \subset \mathcal{P}$. Since $\mathcal{L}_{0}$ is dense in $\mathcal{L}$ and since $\hat{T}$ is a bounded operator, it follows that $\hat{T} \mathcal{L} \subset \mathcal{P}$.
(b) For each $f \in \mathcal{Q}$, it is easy to see that $\tilde{T} J f=J T f$. Combining this with (4.4), we have $\tilde{T} J f=J T f=J J^{*} J f$. Since $J \mathcal{Q}=\mathcal{P}$, this implies $\tilde{T}=J J^{*}$. Since $J: \mathcal{Q} \rightarrow \mathcal{P}$ and $J^{*}: \mathcal{P} \rightarrow \mathcal{Q}$ are invertible, so is $\tilde{T}$.
(c) This follows from (a) and the self-adjointness of $\hat{T}$, which is provided by Lemma 4.5.

In what follows, we write $\zeta_{1}, \ldots, \zeta_{n}$ for the coordinate functions on $\mathbf{C}^{n}$.
Definition 4.7. For $\varphi \in \mathbf{C}\left[\zeta_{1}, \bar{\zeta}_{1}, \ldots, \zeta_{n}, \bar{\zeta}_{n}\right], \hat{M}_{\varphi}$ denotes the operator of multiplication by the function $\varphi$ on $\mathcal{L}$.
Proposition 4.8. For each $j \in\{1, \ldots, n\}, \mathcal{P}$ is an invariant subspace for $\hat{M}_{\zeta_{j}}$.
Proof. Let $f \in \mathcal{Q}$. Then $Q_{\zeta_{j}} f=\zeta_{j} f-g_{j}$ for some $g_{j} \in \mathcal{R}$. It follows from (4.4) that $J g_{j}=0$. Therefore

$$
\begin{equation*}
J Q_{\zeta_{j}} f=J \zeta_{j} f=\zeta_{j} J f=\hat{M}_{\zeta_{j}} J f \tag{4.10}
\end{equation*}
$$

That is, for each $f \in \mathcal{Q}$, we have $\hat{M}_{\zeta_{j}} J f \in J \mathcal{Q}=\mathcal{P}$, which proves the proposition.

Proposition 4.8 makes it possible for us to introduce
Definition 4.9. For each $j \in\{1, \ldots, n\}$, let $M_{\zeta_{j}}$ denote the restriction of the operator $\hat{M}_{\zeta_{j}}$ to the invariant subspace $\mathcal{P}$.

Thus we can restate (4.10) in the form
Corollary 4.10. We have $J Q_{\zeta_{j}}=M_{\zeta_{j}} J$ for every $j \in\{1, \ldots, n\}$.
We end the section with two crucial technical results, whose proofs will have to wait until Section 6.

Proposition 4.11. For every $\varphi \in \mathbf{C}\left[\zeta_{1}, \bar{\zeta}_{1}, \ldots, \zeta_{n}, \bar{\zeta}_{n}\right]$, we have $\left[\hat{M}_{\varphi}, \hat{T}\right] \in \mathcal{C}_{2}$.
Proposition 4.12. For analytic polynomials $q, r \in \mathbf{C}\left[\zeta_{1}, \ldots, \zeta_{n}\right]$, we have

$$
\left[\hat{M}_{\bar{r}},\left[\hat{M}_{q}, \hat{T}\right]\right] P \in \mathcal{C}_{1}
$$

## 5. Operators on $L^{2}(M, d \mu)$

First, we recall the following:
Proposition 5.1. [29, Proposition 11.1] Under the assumption $d=1$, there is a $0<C<$ $\infty$ such that

$$
|\zeta-w| \leq C|1-\langle\zeta, w\rangle|
$$

for all $\zeta, w \in M$.
Lemma 5.2. If $G(z, w)$ is a bounded Borel function on $M \times M$, then the operator

$$
\begin{equation*}
\left(A_{G} \varphi\right)(z)=\int_{M} \frac{G(z, w)}{(1-\langle z, w\rangle)^{3}} \varphi(w) d \mu(w), \quad \varphi \in L^{2}(M, \mu) \tag{5.1}
\end{equation*}
$$

is bounded on $L^{2}(M, \mu)$.
Proof. Consider the function $h(w)=\left(1-|w|^{2}\right)^{-1 / 2}$ on $M$. Recalling (3.1), we have

$$
\begin{aligned}
\int_{M} h(w)\left|\frac{G(z, w)}{(1-\langle z, w\rangle)^{3}}\right| d \mu(w) & \leq\|G\|_{\infty} \int_{M} \frac{\left(1-|w|^{2}\right)^{1 / 2}}{|1-\langle z, w\rangle|^{1+1+(1 / 2)+(1 / 2)}} d v_{M}(w) \\
& \leq C_{1}\|G\|_{\infty} h(z)
\end{aligned}
$$

where the last step is an application of Lemma 2.8 in the case $d=1$. Similarly,

$$
\int_{M} h(z)\left|\frac{G(z, w)}{(1-\langle z, w\rangle)^{3}}\right| d \mu(z) \leq C_{1}\|G\|_{\infty} h(w)
$$

Thus the Schur test gives us $\left\|A_{G}\right\| \leq C_{1}\|G\|_{\infty}$.
On the Hilbert space $L^{2}(M, d \mu)$, we define the operator

$$
(Z \varphi)(z)=\int_{M} \frac{1}{(1-\langle z, w\rangle)^{3}} \varphi(w) d \mu(w), \quad \varphi \in L^{2}(M, d \mu)
$$

Lemma 5.3. If $f$ is a Lipschitz function on $M$, then the commutator $\left[M_{f}, Z\right]$ is a HilbertSchmidt operator on $L^{2}(M, d \mu)$.
Proof. Obviously, the kernel of $\left[M_{f}, Z\right]$ equals $(1-\langle z, w\rangle)^{-3}(f(z)-f(w))$. If $f$ is Lipschitz on $M$, then it follows from Proposition 5.1 that

$$
\iint\left|\frac{f(z)-f(w)}{(1-\langle z, w\rangle)^{3}}\right|^{2} d \mu(w) d \mu(z) \leq C_{1} \iint \frac{1}{|1-\langle z, w\rangle|^{2}} d v_{M}(w) d v_{M}(z)
$$

Note that

$$
\begin{aligned}
\iint \frac{1}{|1-\langle z, w\rangle|^{2}} d v_{M}(w) d v_{M}(z) & \leq \iint \frac{2^{1 / 2}}{|1-\langle z, w\rangle|^{1+1+(1 / 2)}} d v_{M}(w) d v_{M}(z) \\
& \leq C_{2} \int \frac{1}{\left(1-|z|^{2}\right)^{1 / 2}} d v_{M}(z)<\infty
\end{aligned}
$$

where the second $\leq$ is an application of Lemma 2.8. Hence $\left[M_{f}, Z\right] \in \mathcal{C}_{2}$.
By Lemma 5.2, $Z$ is a bounded operator on $L^{2}(M, d \mu)$. Obviously, $Z$ is self-adjoint. What is less obvious is the following:
Lemma 5.4. We have $Z \geq 0$ on $L^{2}(M, d \mu)$.
Proof. To prove this we need the Hilbert space $\mathcal{H}^{(2-n)}$ (see, e.g., [19]). Recall that, since $3=n+1+(2-n), \mathcal{H}^{(2-n)}$ is the Hilbert space of analytic functions on $\mathbf{B}$ which has

$$
\begin{equation*}
K_{w}^{(2-n)}(z)=\frac{1}{(1-\langle z, w\rangle)^{3}}, \quad z, w \in \mathbf{B} \tag{5.2}
\end{equation*}
$$

as its reproducing kernel. That is, the space $\mathcal{H}^{(2-n)}$ has "weight" $2-n$ [19].
For each $0<\alpha<1$,

$$
Y_{\alpha}=\int_{N^{(\alpha)}} K_{w}^{(2-n)} \otimes K_{w}^{(2-n)} d \mu(w)
$$

is obviously a positive operator on $\mathcal{H}^{(2-n)}$ (see (2.13)). For each $0<\alpha<1$, let $I_{\alpha}$ be the operator that maps each $f \in \mathcal{H}^{(2-n)}$ to the function $f \mid N^{(\alpha)}$ in $L^{2}\left(N^{(\alpha)}, d \mu\right)$. Then $\left\|I_{\alpha} f\right\|^{2}=\left\langle Y_{\alpha} f, f\right\rangle$ for every $f \in \mathcal{H}^{(2-n)}$. Therefore $Y_{\alpha}=I_{\alpha}^{*} I_{\alpha}$.

Given an $f \in \mathcal{H}^{(2-n)}$, we have $\left\langle f, K_{w}^{(2-n)}\right\rangle_{2-n}=f(w), w \in \mathbf{B}$. Therefore

$$
\left(I_{\alpha} Y_{\alpha} f\right)(z)=\int_{N^{(\alpha)}} \frac{f(w)}{(1-\langle z, w\rangle)^{3}} d \mu(w) \quad \text { for } \quad z \in N^{(\alpha)} .
$$

On the other hand, by the definitions of $Z$ and $I_{\alpha}$,

$$
\left(M_{\chi_{N}(\alpha)} Z I_{\alpha} f\right)(z)=\int_{N^{(\alpha)}} \frac{f(w)}{(1-\langle z, w\rangle)^{3}} d \mu(w) \quad \text { for } \quad z \in N^{(\alpha)}
$$

Thus under the natural identification of $L^{2}\left(N^{(\alpha)}, d \mu\right)$ with $\left\{\varphi \in L^{2}(M, d \mu): \varphi=0\right.$ on $\left.M \backslash N^{(\alpha)}\right\}$, we have

$$
I_{\alpha} Y_{\alpha}=M_{\chi_{N}(\alpha)} Z I_{\alpha}
$$

Therefore for each analytic polynomial $q \in \mathbf{C}\left[z_{1}, \ldots, z_{n}\right]$,

$$
\left\langle Z I_{\alpha} q, I_{\alpha} q\right\rangle=\left\langle M_{\chi_{N^{(\alpha)}}} Z I_{\alpha} q, I_{\alpha} q\right\rangle=\left\langle I_{\alpha} Y_{\alpha} q, I_{\alpha} q\right\rangle=\left\langle Y_{\alpha} q, I_{\alpha}^{*} I_{\alpha} q\right\rangle=\left\|Y_{\alpha} q\right\|^{2} \geq 0
$$

Since this holds for every $0<\alpha<1$, we have $\langle Z q, q\rangle \geq 0$ for every $q \in \mathbf{C}\left[z_{1}, \ldots, z_{n}\right]$. Since the range of $Z$ is obviously contained in the closure of $\mathbf{C}\left[z_{1}, \ldots, z_{n}\right]$ in $L^{2}(M, d \mu)$ and since $Z$ is self-adjoint, we conclude that $Z$ is positive on $L^{2}(M, d \mu)$.

Definition 5.5. (a) Let $\mathcal{E}$ denote the closure of $\mathbf{C}\left[z_{1}, \ldots, z_{n}\right]$ in $L^{2}(M, d \mu)$.
(b) Let $E$ denote the orthogonal projection from $L^{2}(M, d \mu)$ onto $\mathcal{E}$.

Proposition 5.6. There is a $\gamma>0$ such that the spectrum of $Z$ does not intersect the interval $(0, \gamma)$. Moreover, the range of $Z$ equals $\mathcal{E}$.

The proof of Proposition 5.6 is similar to the work in Section 3, in particular to the proof of Theorem 3.5. For this reason we will present the proof of Proposition 5.6 in Appendix 2.

An immediate consequence of Lemma 5.4 and Proposition 5.6 is that we can write the orthogonal projection $E$ in the form $E=h(Z)$ for some $h \in C_{c}^{\infty}(\mathbf{R})$. Applying the standard smooth functional calculus [9], from Lemma 5.3 we obtain

Corollary 5.7. If $f$ is a Lipschitz function on $M$, then the commutator $\left[M_{f}, E\right]$ is a Hilbert-Schmidt operator on $L^{2}(M, d \mu)$.

We again need the function $\rho(z)=1-|z|^{2}, z \in \mathbf{B}$.
Lemma 5.8. If $r>1 / 2$, then $M_{\rho^{r}} E$ is a Hilbert-Schmidt operator on $L^{2}(M, d \mu)$.
Proof. We know that $Z \geq 0$ from Lemma 5.4. Let $\tilde{Z}$ denote the restriction of $Z$ to its range $\mathcal{E}$. Then Proposition 5.6 implies that $\tilde{Z}$ is invertible and that $E=Z\left(\tilde{Z}^{-1} \oplus 0\right)$. Therefore it suffices to show that $M_{\rho^{r}} Z \in \mathcal{C}_{2}$ for $1 / 2<r<1$.

For each $f \in L^{2}(M, d \mu)$, we have

$$
\left(M_{\rho^{r}} Z f\right)(z)=\int_{M} \frac{\left(1-|z|^{2}\right)^{r}}{(1-\langle z, w\rangle)^{3}} f(w) d \mu(w)
$$

Moreover,

$$
\begin{aligned}
& \iint\left|\frac{\left(1-|z|^{2}\right)^{r}}{(1-\langle z, w\rangle)^{3}}\right|^{2} d \mu(w) d \mu(z) \leq \iint \frac{2^{2+2 r}}{|1-\langle z, w\rangle|^{4-2 r}} d v_{M}(w) d v_{M}(z) \\
& \quad=\iint \frac{2^{2+2 r}}{|1-\langle z, w\rangle|^{1+1+2(1-r)}} d v_{M}(w) d v_{M}(z) \leq \int \frac{C_{1}}{\left(1-|z|^{2}\right)^{2(1-r)}} d v_{M}(z)
\end{aligned}
$$

where the second $\leq$ follows from Lemma 2.8. For $r>1 / 2$, we have $2(1-r)<1$, and the above is finite. Therefore $M_{\rho^{r}} Z \in \mathcal{C}_{2}$.
Lemma 5.9. For all $f, g \in \operatorname{Lip}(M)$ and $0<r<1$, we have $M_{\rho^{-r}}\left[M_{f},\left[M_{g}, Z\right]\right] \in \mathcal{C}_{2}$.
Proof. It suffices to show that

$$
\iint\left\{\rho^{-r}(z)\right\}^{2} \frac{|(f(z)-f(w))(g(z)-g(w))|^{2}}{|1-\langle z, w\rangle|^{6}} d \mu(w) d \mu(z)<\infty
$$

By the Lipschitz condition for $f, g$ and Proposition 5.1, the left-hand side does not exceed

$$
\begin{aligned}
C_{1} \iint\left(1-|z|^{2}\right)^{-2 r} & \frac{|1-\langle z, w\rangle|^{4}}{|1-\langle z, w\rangle|^{6}} d \mu(w) d \mu(z) \leq 2 C_{1} \iint \frac{\left(1-|z|^{2}\right) d v_{M}(w) d v_{M}(z)}{\left(1-|z|^{2}\right)^{2 r}|1-\langle z, w\rangle|} \\
& =2 C_{1} \iint \frac{1}{\left(1-|z|^{2}\right)^{r}} \cdot \frac{\left(1-|z|^{2}\right)^{1-r}}{|1-\langle z, w\rangle|} d v_{M}(w) d v_{M}(z) \\
& \leq 4 C_{1} \iint \frac{d v_{M}(w) d v_{M}(z)}{\left(1-|z|^{2}\right)^{r}\left(1-|w|^{2}\right)^{r}}
\end{aligned}
$$

which is finite because $r<1$. This completes the proof.
Lemma 5.10. For $f, g \in \operatorname{Lip}(M)$, the operators $E\left[M_{f},\left[M_{g}, Z\right]\right]$ and $\left[M_{f},\left[M_{g}, Z\right]\right] E$ are in the trace class $\mathcal{C}_{1}$.
Proof. Take any $1 / 2<r<1$. We have the factorization

$$
E\left[M_{f},\left[M_{g}, Z\right]\right]=E M_{\rho^{r}} \cdot M_{\rho^{-r}}\left[M_{f},\left[M_{g}, Z\right]\right]
$$

Applying Lemmas 5.8 and 5.9, we obtain the membership $E\left[M_{f},\left[M_{g}, Z\right]\right] \in \mathcal{C}_{1}$. Then note that $\left[M_{f},\left[M_{g}, Z\right]\right] E=\left\{E\left[\left[Z, M_{\bar{g}}\right], M_{\bar{f}}\right]\right\}^{*} \in \mathcal{C}_{1}$.
Proposition 5.11. For all $f, g \in \operatorname{Lip}(M)$, the double commutator $\left[M_{f},\left[M_{g}, Z\right]\right]$ is in the trace class.

Proof. Let us write $A \sim_{1} B$ if $A-B \in \mathcal{C}_{1}$. By Lemma 5.10, we have

$$
\left[M_{f},\left[M_{g}, Z\right]\right] \sim_{1}(1-E)\left[M_{f},\left[M_{g}, Z\right]\right] \sim_{1}(1-E)\left[M_{f},\left[M_{g}, Z\right]\right](1-E)
$$

Then note that $(1-E) Z=0$ and $Z(1-E)=0$. Therefore

$$
\begin{aligned}
& (1-E)\left[M_{f},\left[M_{g}, Z\right]\right](1-E)=-(1-E) M_{f} Z M_{g}(1-E)-(1-E) M_{g} Z M_{f}(1-E) \\
& \quad=-\left[1-E, M_{f}\right] Z\left[M_{g}, 1-E\right]-\left[1-E, M_{g}\right] Z\left[M_{f}, 1-E\right] \\
& \quad=\left[M_{f}, E\right] Z\left[M_{g}, E\right]+\left[M_{g}, E\right] Z\left[M_{f}, E\right]
\end{aligned}
$$

which, according to Corollary 5.7, is in the trace class.
On the Hilbert space $L^{2}(M, d \mu)$, we also define the operator

$$
(\Lambda \varphi)(z)=\int_{M} \frac{1}{(1-\langle z, w\rangle)^{2}} \varphi(w) d \mu(w), \quad \varphi \in L^{2}(M, d \mu)
$$

Lemma 5.12. We have $\Lambda \in \mathcal{C}_{2}$. Moreover, $E\left[M_{f}, \Lambda\right] \in \mathcal{C}_{1}$ and $\left[M_{f}, \Lambda\right] E \in \mathcal{C}_{1}$ for every $f \in \operatorname{Lip}(M)$.

Proof. We have

$$
\begin{aligned}
\iint \frac{1}{|1-\langle z, w\rangle|^{4}} d \mu(w) d \mu(z) & \leq \iint \frac{2^{3 / 2}}{|1-\langle z, w\rangle|^{1+1+(1 / 2)}} d v_{M}(w) d v_{M}(z) \\
& \leq \int \frac{C_{1}}{\left(1-|z|^{2}\right)^{1 / 2}} d v_{M}(z)<\infty
\end{aligned}
$$

where the second $\leq$ is an application of Lemma 2.8. Therefore $\Lambda \in \mathcal{C}_{2}$. Similar to what we saw in the proof of Lemma 5.9, we have

$$
M_{\rho^{-r}}\left[M_{f}, \Lambda\right] \in \mathcal{C}_{2}
$$

for $f \in \operatorname{Lip}(M)$ and $1 / 2<r<1$. Combining this membership with the factorization

$$
E\left[M_{f}, \Lambda\right]=E M_{\rho^{r}} \cdot M_{\rho^{-r}}\left[M_{f}, \Lambda\right]
$$

and with Lemma 5.8, we obtain the membership $E\left[M_{f}, \Lambda\right] \in \mathcal{C}_{1}$. By the relation $\left[M_{f}, \Lambda\right] E$ $=\left\{E\left[\Lambda, M_{\bar{f}}\right]\right\}^{*}$, we also have $\left[M_{f}, \Lambda\right] E \in \mathcal{C}_{1}$.

Proposition 5.13. For every $f \in \operatorname{Lip}(M)$ we have $\left[M_{f}, \Lambda\right] \in \mathcal{C}_{1}$.
Proof. Note that $(1-E) \Lambda=0$ and $\Lambda(1-E)=0$. Applying Lemma 5.12, we have

$$
\left[M_{f}, \Lambda\right] \sim_{1}(1-E)\left[M_{f}, \Lambda\right] \sim_{1}(1-E)\left[M_{f}, \Lambda\right](1-E)=0,
$$

i.e., $\left[M_{f}, \Lambda\right] \in \mathcal{C}_{1}$.

Lemma 5.14. Write $\|\cdot\|_{2-n}$ for the norm on the reproducing-kernel Hilbert space $\mathcal{H}^{(2-n)}$. There is a $0<C<\infty$ such that

$$
\int_{M}|f(w)|^{2} d \mu(w) \leq C\|f\|_{2-n}^{2}
$$

for every $f \in \mathcal{H}^{(2-n)}$. In other words, $d \mu$ is a Carleson measure for $\mathcal{H}^{(2-n)}$.
Proof. This is an easier version of Proposition 3.2. We pick a natural number $m \geq 2$ such that $2 m-n \geq 0$. By Lemma 2.14, we have

$$
\begin{aligned}
\int_{M}|f(w)|^{2} d \mu(w) & \leq C_{1} \sum_{|\alpha|=m-1} \int_{M^{(b)}}\left|\left(\partial^{\alpha} f\right)(w)\right|^{2}\left(1-|w|^{2}\right)^{2 m-2} d \mu(w) \\
& +C_{1} \sum_{0 \leq|\beta| \leq m-2} \int_{N^{(a)}}\left|\left(\partial^{\beta} f\right)(w)\right|^{2} d \mu(w),
\end{aligned}
$$

$f \in \mathcal{H}^{(2-n)}$. Recalling (2.13), the second term is dominated by $C_{2}\|f\|_{2-n}^{2}$. It follows from Lemma 2.12 that for each $\alpha \in \mathbf{Z}_{+}^{n}$ with $|\alpha|=m-1$,

$$
\begin{aligned}
& \int_{M^{(b)}}\left|\left(\partial^{\alpha} f\right)(w)\right|^{2}\left(1-|w|^{2}\right)^{2 m-2} d \mu(w) \\
& \quad \leq C_{3} \int_{\mathbf{B}}\left|\left(\partial^{\alpha} f\right)(w)\right|^{2}\left(1-|w|^{2}\right)^{2 m-n} d v(w) \leq C_{4}\|f\|_{2-n}^{2}
\end{aligned}
$$

where the second $\leq$ follows from Lemmas 2.15 and 2.16.

## 6. Proofs of Propositions 4.11 and 4.12

Let $\mathcal{I}: \mathcal{L} \rightarrow L^{2}\left(M^{\left(t_{0}\right)}, d \mu\right)$ be the natural embedding. Obviously, $\|\mathcal{I}\| \leq 1$.
As usual, we identify $L^{2}\left(M^{\left(t_{0}\right)}, d \mu\right)$ with the subspace $\left\{f \in L^{2}(M, d \mu): f=0\right.$ on $\left.M \backslash M^{\left(t_{0}\right)}\right\}$ of $L^{2}(M, d \mu)$.

Lemma 6.1. If $r<1$, then the operator $M_{\rho^{-r}} \mathcal{I}$ is bounded.
Proof. Denote $\tau=1-2 r$. If $r<1$, then $\tau>-1$, and therefore we can apply Lemma 2.13 to the measure $d \nu_{\tau}$. By that lemma, there are $0<a<b<t_{0}$ and $0<C<\infty$ such that for each $f \in \mathcal{L}_{0}$, we have

$$
\begin{aligned}
& \int_{M^{\left(t_{0}\right)}}\left|\rho^{-r}(w) f(w)\right|^{2} d \mu(w)=\int_{M^{\left(t_{0}\right)}}|f(w)|^{2} d \nu_{\tau}(w) \\
& \quad \leq C \int_{M^{(b)}}\left|\left(\nabla_{\mathcal{M}} f\right)(w)\right|^{2}\left(1-|w|^{2}\right)^{2} d \nu_{\tau}(w)+C \int_{N^{(a)} \cap M^{\left(t_{0}\right)}}|f(w)|^{2} d \nu_{\tau}(w)
\end{aligned}
$$

It follows from (2.2) that if $f$ is real-valued, then $\left|\left(\nabla_{\mathcal{M}} f\right)(w)\right| \leq(2 / \gamma)\left|\left(\partial_{p_{w}} f\right)(w)\right|, w \in$ $M^{\left(t_{0}\right)}$. Substituting this in the above and noting that $2+\tau>1$, we have

$$
\begin{aligned}
\int_{M^{\left(t_{0}\right)}} \mid \rho^{-r}(w) & \left.f(w)\right|^{2} d \mu(w) \\
\leq & (2 / \gamma)^{2} C \int_{M^{(b)}}\left|\left(\partial_{p_{w}} f\right)(w)\right|^{2} d \mu(w)+a^{-2} C \int_{N^{(a)} \cap M^{\left(t_{0}\right)}}|f(w)|^{2} d \mu(w)
\end{aligned}
$$

for real-valued $f \in \mathcal{L}_{0}$. Hence there is a $0<C_{1}<\infty$ such that $\left\|M_{\rho^{-r}} \mathcal{I} f\right\| \leq C_{1}\|f\|_{\#}$ for real-valued $f \in \mathcal{L}$.

Now we use the fact that the norm $\|\cdot\|_{\#}$ has the symmetry $\|\bar{g}\|_{\#}=\|g\|_{\#}$ for every $g \in \mathcal{L}$. Thus if $g=f_{1}+i f_{2}$, where $f_{1}$ and $f_{2}$ are real-valued, then $\left\|f_{1}\right\|_{\#} \leq\|g\|_{\#}$ and $\left\|f_{2}\right\|_{\#} \leq\|g\|_{\#}$. Combining this fact with the conclusion of the preceding paragraph, we see that $\left\|M_{\rho^{-r}} \mathcal{I} g\right\| \leq 2 C_{1}\|g\|_{\#}$ for every $g \in \mathcal{L}$.
Lemma 6.2. If $\varphi \in L^{\infty}\left(M^{\left(t_{0}\right)}, d \mu\right)$, then $E M_{\varphi} \mathcal{I}$ is a Hilbert-Schmidt operator.
Proof. Take any $1 / 2<r<1$ and factor $E M_{\varphi} \mathcal{I}$ in the form

$$
E M_{\varphi} \mathcal{I}=E M_{\rho^{r}} \cdot M_{\varphi} \cdot M_{\rho^{-r}} \mathcal{I}
$$

Since $E M_{\rho^{r}}=\left(M_{\rho^{r}} E\right)^{*}$, it follows from Lemmas 5.8 and 6.1 that $E M_{\varphi} \mathcal{I}$ is a HilbertSchmidt operator.

We will write $I$ for the restriction of $\mathcal{I}$ to the subspace $\mathcal{P}$.
Lemma 6.3. The embedding $I: \mathcal{P} \rightarrow L^{2}\left(M^{\left(t_{0}\right)}, d \mu\right)$ is a Hilbert-Schmidt operator.
Proof. Recall from Section 4 that the operator $J: \mathcal{Q} \rightarrow \mathcal{P}$ is invertible. Thus it suffices to show that $I J: \mathcal{Q} \rightarrow L^{2}\left(M^{\left(t_{0}\right)}, d \mu\right)$ is a Hilbert-Schmidt operator. By Lemma 2.13, there are $0<a<b<1$ and $0<C<\infty$ such that

$$
\int_{M}|f(w)|^{2} d \mu(w) \leq C \int_{M^{(b)}}\left|\left(\nabla_{\mathcal{M}} f\right)(w)\right|^{2}\left(1-|w|^{2}\right)^{2} d \mu(w)+C \int_{N^{(a)}}|f(w)|^{2} d \mu(w)
$$

for every $f \in H_{n}^{2}$. Therefore

$$
\begin{align*}
& \int_{M}|f(w)|^{2} d \mu(w) \\
& 6.1) \quad \leq 2 C \sum_{j=1}^{n} \int_{M}\left|\left(\partial_{j} f\right)(w)\right|^{2}\left(1-|w|^{2}\right)^{2} d \mu(w)+\frac{C}{a^{2}} \int_{M}|f(w)|^{2}\left(1-|w|^{2}\right)^{2} d \mu(w) \tag{6.1}
\end{align*}
$$

Write $\|\cdot\|_{2-n}$ for the norm on $\mathcal{H}^{(2-n)}$. Then it is well known that $\left\|\partial_{j} g\right\|_{2-n} \leq C_{1}\|g\|$ for all $g \in H_{n}^{2}$ and $1 \leq j \leq n$. Moreover, $\|g\|_{2-n} \leq\|g\|$ for $g \in H_{n}^{2}$. Combining these bounds with Lemma 5.14, we see that

$$
A f=\left(\partial_{1} f, \ldots, \partial_{n} f, f\right), \quad f \in \mathcal{Q}
$$

is a bounded operator that maps $\mathcal{Q}$ into $\mathcal{E}^{[n+1]}$, the orthogonal sum of $n+1$ copies of $\mathcal{E}$. Let $\left\{M_{\rho} E\right\}^{[n+1]}$ denote the orthogonal sum of $n+1$ copies of $M_{\rho} E$. By (6.1), there is a bounded operator $B:\left\{L^{2}(M, d \mu)\right\}^{[n+1]} \rightarrow L^{2}\left(M^{\left(t_{0}\right)}, d \mu\right)$ such that

$$
I J=B\left\{M_{\rho} E\right\}^{[n+1]} A .
$$

By Lemma 5.8, $\left\{M_{\rho} E\right\}^{[n+1]}$ is a Hilbert-Schmidt operator. Therefore so is $I J: \mathcal{Q} \rightarrow$ $L^{2}\left(M^{\left(t_{0}\right)}, d \mu\right)$. This completes the proof.

For each $f \in \mathcal{L}_{0}$, we define

$$
(\mathcal{D} f)(z)=\left(\partial_{p_{z}} f\right)(z) \quad \text { and } \quad(\overline{\mathcal{D}} f)(z)=\left(\overline{\partial_{p_{z}}} f\right)(z)
$$

We consider $\mathcal{D}$ and $\overline{\mathcal{D}}$ as operators from $\mathcal{L}_{0}$ into $L^{2}\left(M^{\left(t_{0}\right)}, d \mu\right)$. By the definition of $\|\cdot\|_{\#}$, $\mathcal{D}$ and $\overline{\mathcal{D}}$ are contractions on $\mathcal{L}_{0}$. Thus $\mathcal{D}$ and $\overline{\mathcal{D}}$ naturally extend to contractions form $\mathcal{L}$ to $L^{2}\left(M^{\left(t_{0}\right)}, d \mu\right)$.

We will write $D$ for the restriction of $\mathcal{D}$ to the subspace $\mathcal{P}$. Note that the restriction of $\overline{\mathcal{D}}$ to $\mathcal{P}$ is 0 .

Notation 6.4. For the rest of the paper, $\psi_{i}(\zeta)$ denotes the $i$-the component of the vector $p_{\zeta}, \zeta \in \mathcal{M}$ and $1 \leq i \leq n$.

Lemma 6.5. If $\varphi \in L^{\infty}\left(M^{\left(t_{0}\right)}, d \mu\right)$, then $\mathcal{I}^{*} M_{\varphi} \mathcal{D} \hat{T}_{1} \in \mathcal{C}_{2}$ and $\mathcal{I}^{*} M_{\varphi} \mathcal{D} \hat{T}_{2} \in \mathcal{C}_{2}$.
Proof. For $f \in \mathcal{L}_{0}$, straightforward differentiation gives us

$$
\left(\mathcal{D} \hat{T}_{1} f\right)(z)=2(A f)(z)+(B f)(z)
$$

where

$$
\begin{aligned}
& (A f)(z)=\int_{M^{\left(t_{0}\right)}} \frac{\left\langle p_{z}, w\right\rangle\left\langle z, p_{w}\right\rangle}{(1-\langle z, w\rangle)^{3}}\left(\partial_{p_{w}} f\right)(w) d \mu(w) \quad \text { and } \\
& (B f)(z)=\int_{M^{\left(t_{0}\right)}} \frac{\left\langle p_{z}, p_{w}\right\rangle}{(1-\langle z, w\rangle)^{2}}\left(\partial_{p_{w}} f\right)(w) d \mu(w)
\end{aligned}
$$

Thus

$$
A=\sum_{j=1}^{n} \sum_{i=1}^{n} M_{\chi_{M^{\left(t_{0}\right)}}} M_{\psi_{j}} M_{\zeta_{i}} Z M_{\bar{\zeta}_{j}} M_{\bar{\psi}_{i}} \mathcal{D} \quad \text { and } \quad B=\sum_{i=1}^{n} M_{\chi_{M^{\left(t_{0}\right)}}} M_{\psi_{i}} \Lambda M_{\bar{\psi}_{i}} \mathcal{D} .
$$

Since $Z=E Z$, for $1 \leq i, j \leq n$ we have

$$
\mathcal{I}^{*} M_{\varphi} M_{\chi_{M^{\left(t_{0}\right)}}} M_{\psi_{j}} M_{\zeta_{i}} Z=\mathcal{I}^{*} M_{\varphi \chi_{M}\left(t_{0}\right)} \psi_{j} \zeta_{i} E Z=\left(E M_{\bar{\varphi} \chi_{M}\left(t_{0}\right)} \bar{\psi}_{j} \bar{\zeta}_{i} \mathcal{I}\right)^{*} Z
$$

Thus from Lemma 6.2 we see that $\mathcal{I}^{*} M_{\varphi} A$ is a Hilbert-Schmidt operator. On the other hand, it follows from Lemma 5.12 that $B$ is a Hilbert-Schmidt operator. Therefore $\mathcal{I}^{*} M_{\varphi} \mathcal{D} \hat{T}_{1}=\mathcal{I}^{*} M_{\varphi}(2 A+B)$ is a Hilbert-Schmidt operator.

Straightforward differentiation gives us

$$
\left(\mathcal{D} \hat{T}_{2} f\right)(z)=\int_{M^{\left(t_{0}\right)}} \frac{\left\langle p_{z}, w\right\rangle}{(1-\langle z, w\rangle)^{2}} f(w) d \mu(w)
$$

$f \in \mathcal{L}_{0}$. That is,

$$
\mathcal{D} \hat{T}_{2}=\sum_{i=1}^{n} M_{\psi_{i}} \Lambda M_{\bar{\zeta}_{i}} \mathcal{I}
$$

By Lemma $5.12, \mathcal{D} \hat{T}_{2}$ is a Hilbert-Schmidt operator, and so is $\mathcal{I}^{*} M_{\varphi} \mathcal{D} \hat{T}_{2}$. This completes the proof.

Since $\mathcal{L}$ is not an $L^{2}$-space, the adjoint of $\hat{M}_{\zeta_{j}}$ is not $\hat{M}_{\bar{\zeta}_{j}}, 1 \leq j \leq n$. But using the operators $\mathcal{I}, \mathcal{D}$ and $\overline{\mathcal{D}}$, we can give a formula for $\hat{M}_{\zeta_{j}}^{*}$ :

Proposition 6.6. For each $1 \leq j \leq n$, we have $\hat{M}_{\zeta_{j}}^{*}=\hat{M}_{\bar{\zeta}_{j}}+\mathcal{I}^{*} M_{\bar{\psi}_{j}} \mathcal{D}-\overline{\mathcal{D}}^{*} M_{\bar{\psi}_{j}} \mathcal{I}$.

Proof. For all $f, g \in \mathcal{L}_{0}$ and $j \in\{1, \ldots, n\}$, we have

$$
\begin{aligned}
&\left\langle f, \hat{M}_{\zeta_{j}}^{*} g\right\rangle_{\#}=\left\langle\hat{M}_{\zeta_{j}} f, g\right\rangle_{\#}=\int_{M^{\left(t_{0}\right)}}\left(\partial_{p_{w}} \hat{M}_{\zeta_{j}} f\right)(w) \overline{\left(\partial_{p_{w}} g\right)(w)} d \mu(w) \\
&+\int_{M^{\left(t_{0}\right)}}\left(\overline{\partial_{p_{w}}} \hat{M}_{\zeta_{j}} f\right)(w) \overline{\left(\overline{\partial_{p_{w}}} g\right)(w)} d \mu(w)+\int_{M^{\left(t_{0}\right)}} w_{j} f(w) \overline{g(w)} d \mu(w) \\
&= \int_{M^{\left(t_{0}\right)}} w_{j}\left(\partial_{p_{w}} f\right)(w) \overline{\left(\partial_{p_{w}} g\right)(w)} d \mu(w)+\int_{M^{\left(t_{0}\right)}} w_{j}\left(\overline{\partial_{p_{w}}} f\right)(w) \overline{\left(\overline{\partial_{p_{w}}} g\right)(w)} d \mu(w) \\
&+\int_{M^{\left(t_{0}\right)}} \psi_{j}(w) f(w) \overline{\left(\partial_{p_{w}} g\right)(w)} d \mu(w)+\int_{M^{\left(t_{0}\right)}} w_{j} f(w) \overline{g(w)} d \mu(w) \\
&= \int_{M^{\left(t_{0}\right)}}\left(\partial_{p_{w}} f\right)(w) \overline{\left(\partial_{p_{w}} \hat{M}_{\bar{\zeta}_{j}} g\right)(w)} d \mu(w)+\int_{M^{\left(t_{0}\right)}}\left(\overline{\partial_{p_{w}}} f\right)(w) \overline{\left(\overline{\partial_{p_{w}}} \hat{M}_{\bar{\zeta}_{j}} g\right)(w)} d \mu(w) \\
&+\int_{M^{\left(t_{0}\right)}} f(w) \overline{\left(\hat{M}_{\bar{\zeta}_{j}} g\right)(w)} d \mu(w) \\
&+\int_{M^{\left(t_{0}\right)}} \psi_{j}(w) f(w) \overline{\left(\partial_{p_{w}} g\right)(w)} d \mu(w)-\int_{M^{\left(t_{0}\right)}} \psi_{j}(w)\left(\overline{\partial_{p_{w}}} f\right)(w) \overline{g(w)} d \mu(w) \\
&=\left\langle f, \hat{M}_{\bar{\zeta}_{j}} g\right\rangle_{\#}+\left\langle M_{\psi_{j}} \mathcal{I} f, \mathcal{D} g\right\rangle-\left\langle M_{\psi_{j}} \overline{\mathcal{D}} f, \mathcal{I} g\right\rangle \\
&=\left\langle f, \hat{M}_{\bar{\zeta}_{j}} g\right\rangle_{\#}+\left\langle f, \mathcal{I}^{*} M_{\bar{\psi}_{j}} \mathcal{D} g\right\rangle_{\#}-\left\langle f, \overline{\mathcal{D}}^{*} M_{\bar{\psi}_{j}} \mathcal{I} g\right\rangle_{\#} .
\end{aligned}
$$

This completes the proof.
Lemma 6.7. If $\varphi \in L^{\infty}\left(M^{\left(t_{0}\right)}, d \mu\right)$, then $\left[\mathcal{I}^{*} M_{\varphi} \mathcal{D}, \hat{T}\right] \in \mathcal{C}_{2}$ and $\left[\overline{\mathcal{D}}^{*} M_{\varphi} \mathcal{I}, \hat{T}\right] \in \mathcal{C}_{2}$.
Proof. We have $\mathcal{I}^{*} M_{\varphi} \mathcal{D} \hat{T} \in \mathcal{C}_{2}$ by Lemma 6.5. On the other hand, by Lemma 6.3, $\hat{T} \mathcal{I}^{*}=$ $(\mathcal{I} \hat{T})^{*}=(\mathcal{I} P \hat{T})^{*} \in \mathcal{C}_{2}$. Therefore $\left[\mathcal{I}^{*} M_{\varphi} \mathcal{D}, \hat{T}\right] \in \mathcal{C}_{2}$.

Similarly, $\mathcal{I} \hat{T}=\mathcal{I} P \hat{T} \in \mathcal{C}_{2}$ by Lemma 6.3. Then note that $\hat{T} \overline{\mathcal{D}}^{*}=(\overline{\mathcal{D}} \hat{T})^{*}=0$. Hence $\left[\overline{\mathcal{D}}^{*} M_{\varphi} \mathcal{I}, \hat{T}\right] \in \mathcal{C}_{2}$.
Lemma 6.8. For every $1 \leq j \leq n$, we have $\left[\hat{M}_{\bar{\zeta}_{j}}, \hat{T}\right](1-P) \in \mathcal{C}_{2}$.
Proof. It follows from Proposition 6.6 that

$$
\left[\hat{M}_{\bar{\zeta}_{j}}, \hat{T}\right](1-P)=\left[\hat{M}_{\zeta_{j}}^{*}, \hat{T}\right](1-P)-\left[\mathcal{I}^{*} M_{\bar{\psi}_{j}} \mathcal{D}, \hat{T}\right](1-P)+\left[\overline{\mathcal{D}}^{*} M_{\bar{\psi}_{j}} \mathcal{I}, \hat{T}\right](1-P)
$$

By Lemma 6.7, the last two terms on the right-hand side are in $\mathcal{C}_{2}$. Then note that

$$
\left[\hat{M}_{\zeta_{j}}^{*}, \hat{T}\right](1-P)=\left\{(1-P)\left[\hat{T}, \hat{M}_{\zeta_{j}}\right]\right\}^{*}=0
$$

This completes the proof.
Lemma 6.9. Let $A \in \mathcal{B}(\mathcal{L})$. For any $1 \leq p<\infty$, if the operators $\mathcal{D} A, \overline{\mathcal{D}} A, \mathcal{I} A: \mathcal{L} \rightarrow$ $L^{2}\left(M^{\left(t_{0}\right)}, d \mu\right)$ are in the Schatten p-class, then $A$ is in the Schatten p-class.

Proof. By the definition of $\|\cdot\|_{\#}$, the formula

$$
V f=\mathcal{D} f \oplus \overline{\mathcal{D}} f \oplus \mathcal{I} f
$$

$f \in \mathcal{L}$, defines an isometry that maps $\mathcal{L}$ into $L^{2}\left(M^{\left(t_{0}\right)}, d \mu\right) \oplus L^{2}\left(M^{\left(t_{0}\right)}, d \mu\right) \oplus L^{2}\left(M^{\left(t_{0}\right)}, d \mu\right)$. Thus we can factor any operator $A$ on $\mathcal{L}$ in the form

$$
A=V^{*}(\mathcal{D} A \oplus \overline{\mathcal{D}} A \oplus \mathcal{I} A) W
$$

where $W: \mathcal{L} \rightarrow \mathcal{L} \oplus \mathcal{L} \oplus \mathcal{L}$ is the operator defined by the formula $W f=f \oplus f \oplus f, f \in \mathcal{L}$. Obviously, the desire conclusion follows from this factorization.
Lemma 6.10. For each $\varphi \in \mathbf{C}\left[\zeta_{1}, \bar{\zeta}_{1}, \ldots, \zeta_{n}, \bar{\zeta}_{n}\right]$, the operator $\hat{T}_{2} \hat{M}_{\varphi} P$ is in the trace class. Proof. Let $\varphi \in \mathbf{C}\left[\zeta_{1}, \bar{\zeta}_{1}, \ldots, \zeta_{n}, \bar{\zeta}_{n}\right]$. By Lemma 6.9, it suffices to show that the operators $\mathcal{D} \hat{T}_{2} \hat{M}_{\varphi}: \mathcal{P} \rightarrow L^{2}\left(M^{\left(t_{0}\right)}, d \mu\right)$ and $\mathcal{I} \hat{T}_{2} \hat{M}_{\varphi}: \mathcal{P} \rightarrow L^{2}\left(M^{\left(t_{0}\right)}, d \mu\right)$ are in the trace class.

For each $f \in \mathcal{P}$, we have

$$
\left(\mathcal{D} \hat{T}_{2} \hat{M}_{\varphi} f\right)(z)=\int_{M^{\left(t_{0}\right)}} \frac{\left\langle p_{z}, w\right\rangle}{(1-\langle z, w\rangle)^{2}} \varphi(w) f(w) d \mu(w)
$$

This obviously translates to the operator identity

$$
\mathcal{D} \hat{T}_{2} \hat{M}_{\varphi} P=\sum_{i=1}^{n} M_{\chi_{M^{\left(t_{0}\right)}}} M_{\psi_{i}} \Lambda M_{\bar{\zeta}_{i}} M_{\varphi} I
$$

Thus it follows from Lemmas 5.12 and 6.3 that $\mathcal{D} \hat{T}_{2} \hat{M}_{\varphi} P$ is in the trace class. On the other hand, since $(1-\langle z, w\rangle)^{-1}=(1-\langle z, w\rangle)^{-2} \cdot(1-\langle z, w\rangle)$, we have

$$
\mathcal{I} \hat{T}_{2} \hat{M}_{\varphi} P=M_{\chi_{M^{\left(t_{0}\right)}}} \Lambda M_{\varphi} I-\sum_{i=1}^{n} M_{\chi_{M^{\left(t_{0}\right)}}} M_{\zeta_{i}} \Lambda M_{\bar{\zeta}_{i}} M_{\varphi} I
$$

By Lemmas 5.12 and $6.3, \mathcal{I} \hat{T}_{2} \hat{M}_{\varphi} P$ is also in the trace class.
Proof of Proposition 4.11. We begin with a sequence of reductions of our task.
By the linearity of commutators, to prove that $\left[\hat{M}_{\varphi}, \hat{T}\right] \in \mathcal{C}_{2}$ for an arbitrary $\varphi \in$ $\mathbf{C}\left[\zeta_{1}, \bar{\zeta}_{1}, \ldots, \zeta_{n}, \bar{\zeta}_{n}\right]$, it suffices to show that $\left[\hat{M}_{\zeta^{\alpha}} \bar{\zeta}^{\beta}, \hat{T}\right] \in \mathcal{C}_{2}$ for all $\alpha, \beta \in \mathbf{Z}_{+}^{n}$. Since

$$
\left[A_{1} A_{2} \cdots A_{\nu}, \hat{T}\right]=\left[A_{1}, \hat{T}\right] A_{2} \cdots A_{\nu}+A_{1}\left[A_{2}, \hat{T}\right] A_{3} \cdots A_{\nu}+\cdots+A_{1} \cdots A_{\nu-1}\left[A_{\nu}, \hat{T}\right]
$$

our task is reduced to the proof of the fact that $\left[\hat{M}_{\zeta_{k}}, \hat{T}\right] \in \mathcal{C}_{2}$ and $\left[\hat{M}_{\bar{\zeta}_{k}}, \hat{T}\right] \in \mathcal{C}_{2}$ for each $1 \leq k \leq n$. By Proposition 6.6, we have

$$
\left[\hat{T}, \hat{M}_{\zeta_{k}}\right]^{*}=\left[\hat{M}_{\zeta_{k}}^{*}, \hat{T}\right]=\left[\hat{M}_{\bar{\zeta}_{k}}, \hat{T}\right]+\left[\mathcal{I}^{*} M_{\bar{\psi}_{k}} \mathcal{D}, \hat{T}\right]-\left[\overline{\mathcal{D}}^{*} M_{\bar{\psi}_{k}} \mathcal{I}, \hat{T}\right] .
$$

Lemma 6.7 tells us that the last two terms on the right-hand side are in $\mathcal{C}_{2}$. Therefore, to prove Proposition 4.11, it suffices to show that $\left[\hat{M}_{\bar{\zeta}_{k}}, \hat{T}\right] \in \mathcal{C}_{2}$ for every $1 \leq k \leq n$. Lemma 6.8 further reduces this task to the proof of the membership $\left[\hat{M}_{\bar{\zeta}_{k}}, \hat{T}\right] P \in \mathcal{C}_{2}, 1 \leq k \leq n$.

By Lemma 6.10 , we actually have $\left[\hat{M}_{\bar{\zeta}_{k}}, \hat{T}_{2}\right] P \in \mathcal{C}_{1}, 1 \leq k \leq n$. Since $\hat{T}=\hat{T}_{1}+\hat{T}_{2}$, we only need to show that $\left[\hat{M}_{\bar{\zeta}_{k}}, \hat{T}_{1}\right] \in \mathcal{C}_{2}$ for every $1 \leq k \leq n$.

By Lemma 6.9, to prove that $\left[\hat{M}_{\bar{\zeta}_{k}}, \hat{T}_{1}\right] \in \mathcal{C}_{2}$, it suffices to show that $\mathcal{D}\left[\hat{M}_{\bar{\zeta}_{k}}, \hat{T}_{1}\right]$, $\overline{\mathcal{D}}\left[\hat{M}_{\bar{\zeta}_{k}}, \hat{T}_{1}\right]$ and $\mathcal{I}\left[\hat{M}_{\bar{\zeta}_{k}}, \hat{T}_{1}\right]$ are Hilbert-Schmidt operators. For $f \in \mathcal{L}_{0}$, we have

$$
\left(\left[\hat{M}_{\bar{\zeta}_{k}}, \hat{T}_{1}\right] f\right)(z)=\int_{M^{\left(t_{0}\right)}} \frac{\bar{z}_{k}-\bar{w}_{k}}{(1-\langle z, w\rangle)^{2}}\left\langle z, p_{w}\right\rangle\left(\partial_{p_{w}} f\right)(w) d \mu(w)
$$

By straightforward differentiation,

$$
\left(\mathcal{D}\left[\hat{M}_{\bar{\zeta}_{k}}, \hat{T}_{1}\right] f\right)(z)=\left(Y_{1} f\right)(z)+\left(Y_{2} f\right)(z)
$$

where

$$
\begin{aligned}
& \left(Y_{1} f\right)(z)=2 \int_{M^{\left(t_{0}\right)}} \frac{\bar{z}_{k}-\bar{w}_{k}}{(1-\langle z, w\rangle)^{3}}\left\langle p_{z}, w\right\rangle\left\langle z, p_{w}\right\rangle\left(\partial_{p_{w}} f\right)(w) d \mu(w) \quad \text { and } \\
& \left(Y_{2} f\right)(z)=\int_{M^{\left(t_{0}\right)}} \frac{\bar{z}_{k}-\bar{w}_{k}}{(1-\langle z, w\rangle)^{2}}\left\langle p_{z}, p_{w}\right\rangle\left(\partial_{p_{w}} f\right)(w) d \mu(w)
\end{aligned}
$$

It is easy to see that

$$
\begin{aligned}
Y_{1} & =2 \sum_{i=1}^{n} \sum_{j=1}^{n} M_{\chi_{M^{\left(t_{0}\right)}}} M_{\psi_{i}} M_{\zeta_{j}}\left[M_{\bar{\zeta}_{k}}, Z\right] M_{\bar{\zeta}_{i}} M_{\bar{\psi}_{j}} \mathcal{D} \quad \text { and } \\
Y_{2} & =\sum_{j=1}^{n} M_{\chi_{M^{\left(t_{0}\right)}}} M_{\psi_{j}}\left[M_{\bar{\zeta}_{k}}, \Lambda\right] M_{\bar{\psi}_{j}} \mathcal{D} .
\end{aligned}
$$

By Lemma 5.3, $Y_{1}$ is a Hilbert-Schmidt operator. By Lemma 5.12, $Y_{2}$ is also a HilbertSchmidt operator. Hence $\mathcal{D}\left[\hat{M}_{\bar{\zeta}_{k}}, \hat{T}_{1}\right]=Y_{1}+Y_{2}$ is a Hilbert-Schmidt operator.

Again by differentiation,

$$
\left(\overline{\mathcal{D}}\left[\hat{M}_{\bar{\zeta}_{k}}, \hat{T}_{1}\right] f\right)(z)=\int_{M^{\left(t_{0}\right)}} \frac{\bar{\psi}_{k}(z)}{(1-\langle z, w\rangle)^{2}}\left\langle z, p_{w}\right\rangle\left(\partial_{p_{w}} f\right)(w) d \mu(w)
$$

for $f \in \mathcal{L}_{0}$. That is,

$$
\overline{\mathcal{D}}\left[\hat{M}_{\bar{\zeta}_{k}}, \hat{T}_{1}\right]=\sum_{j=1}^{n} M_{\chi_{M^{\left(t_{0}\right)}}} M_{\bar{\psi}_{k}} M_{\zeta_{j}} \Lambda M_{\bar{\psi}_{j}} \mathcal{D} .
$$

Thus it follows from Lemma 5.12 that $\overline{\mathcal{D}}\left[\hat{M}_{\bar{\zeta}_{k}}, \hat{T}_{1}\right]$ is a Hilbert-Schmidt operator.
It is easy to see that

$$
\mathcal{I}\left[\hat{M}_{\bar{\zeta}_{k}}, \hat{T}_{1}\right]=\sum_{i=1}^{n} M_{\chi_{M^{\left(t_{0}\right)}}} M_{\zeta_{i}}\left[M_{\bar{\zeta}_{k}}, \Lambda\right] M_{\bar{\psi}_{i}} \mathcal{D}
$$

which is in $\mathcal{C}_{2}$ according to Lemma 5.12. Combining this with the results of the previous two paragraphs, we obtain the membership $\left[\hat{M}_{\bar{\zeta}_{k}}, \hat{T}_{1}\right] \in \mathcal{C}_{2}$. This completes the proof.
Corollary 6.11. For every $\varphi \in \mathbf{C}\left[\zeta_{1}, \bar{\zeta}_{1}, \ldots, \zeta_{n}, \bar{\zeta}_{n}\right]$, we have $\left[\hat{M}_{\varphi}, P\right] \in \mathcal{C}_{2}$.
Proof. By Proposition 4.6, there is an $h \in C_{c}^{\infty}(\mathbf{R})$ such that $P=h(\hat{T})$. Thus the membership $\left[\hat{M}_{\varphi}, P\right] \in \mathcal{C}_{2}$ follows from Proposition 4.11 and the standard smooth functional calculus.

The proof of Proposition 4.11 gives us a taste of what is to come. For Proposition 4.12, because it involves double commutators, the proof will be more tedious for an obvious reason: more terms will have to be examined.

Proof of Proposition 4.12. Let $q, r \in \mathbf{C}\left[\zeta_{1}, \ldots, \zeta_{n}\right]$ be given. By the relation $\hat{T}=\hat{T}_{1}+\hat{T}_{2}$ and Lemma 6.10, it suffices to prove that $\left[\hat{M}_{\bar{r}},\left[\hat{M}_{q}, \hat{T}_{1}\right]\right] P \in \mathcal{C}_{1}$.

For $f \in \mathcal{L}_{0}$, we have

$$
\left(\left[\hat{M}_{q}, \hat{T}_{1}\right] f\right)(z)=\left(X_{1} f\right)(z)-\left(X_{2} f\right)(z)
$$

where

$$
\begin{aligned}
& \left(X_{1} f\right)(z)=\int_{M^{\left(t_{0}\right)}} \frac{q(z)-q(w)}{(1-\langle z, w\rangle)^{2}}\left\langle z, p_{w}\right\rangle\left(\partial_{p_{w}} f\right)(w) d \mu(w) \quad \text { and } \\
& \left(X_{2} f\right)(z)=\int_{M^{\left(t_{0}\right)}} \frac{\left(\partial_{p_{w}} q\right)(w)}{(1-\langle z, w\rangle)^{2}}\left\langle z, p_{w}\right\rangle f(w) d \mu(w)
\end{aligned}
$$

Denote $A_{1}=\left[\hat{M}_{\bar{r}}, X_{1}\right] P$. For $f \in \mathcal{P}$, we have

$$
\left(A_{1} f\right)(z)=\int_{M^{\left(t_{0}\right)}} \frac{(q(z)-q(w))(\bar{r}(z)-\bar{r}(w))}{(1-\langle z, w\rangle)^{2}}\left\langle z, p_{w}\right\rangle\left(\partial_{p_{w}} f\right)(w) d \mu(w)
$$

We again use the operators $\mathcal{D}, \overline{\mathcal{D}}$ and $\mathcal{I}$. By differentiation, we have

$$
\left(\mathcal{D} A_{1} f\right)(z)=\left(Y_{11} f\right)(z)+\left(Y_{12} f\right)(z)+\left(Y_{13} f\right)(z)
$$

where

$$
\begin{aligned}
& \left(Y_{11} f\right)(z)=2 \int_{M^{\left(t_{0}\right)}} \frac{(q(z)-q(w))(\bar{r}(z)-\bar{r}(w))}{(1-\langle z, w\rangle)^{3}}\left\langle p_{z}, w\right\rangle\left\langle z, p_{w}\right\rangle\left(\partial_{p_{w}} f\right)(w) d \mu(w), \\
& \left(Y_{12} f\right)(z)=\int_{M^{\left(t_{0}\right)}} \frac{\left(\partial_{p_{z}} q\right)(z)(\bar{r}(z)-\bar{r}(w))}{(1-\langle z, w\rangle)^{2}}\left\langle z, p_{w}\right\rangle\left(\partial_{p_{w}} f\right)(w) d \mu(w) \quad \text { and } \\
& \left(Y_{13} f\right)(z)=\int_{M^{\left(t_{0}\right)}} \frac{(q(z)-q(w))(\bar{r}(z)-\bar{r}(w))}{(1-\langle z, w\rangle)^{2}}\left\langle p_{z}, p_{w}\right\rangle\left(\partial_{p_{w}} f\right)(w) d \mu(w)
\end{aligned}
$$

$f \in \mathcal{P}$. Denote $\eta(z)=\left(\partial_{p_{z}} q\right)(z)$. We can rewrite the above as the operator identities

$$
\begin{aligned}
& Y_{11}=2 \sum_{i=1}^{n} \sum_{j=1}^{n} M_{\chi_{M^{\left(t_{0}\right)}}} M_{\zeta_{j}} M_{\psi_{i}}\left[M_{\bar{r}},\left[M_{q}, Z\right]\right] M_{\bar{\zeta}_{i}} M_{\bar{\psi}_{j}} D \\
& Y_{12}=\sum_{j=1}^{n} M_{\chi_{M^{\left(t_{0}\right)}}} M_{\eta} M_{\zeta_{j}}\left[M_{\bar{r}}, \Lambda\right] M_{\bar{\psi}_{j}} D \quad \text { and } \\
& Y_{13}=\sum_{j=1}^{n} M_{\chi_{M^{\left(t_{0}\right)}}} M_{\psi_{j}}\left[M_{\bar{r}},\left[M_{q}, \Lambda\right]\right] M_{\bar{\psi}_{j}} D
\end{aligned}
$$

It follows from Proposition 5.11 that $Y_{11} \in \mathcal{C}_{1}$. By Proposition 5.13, we have $Y_{12} \in \mathcal{C}_{1}$ and $Y_{13} \in \mathcal{C}_{1}$. Hence $\mathcal{D} A_{1}$ is in the trace class.

On the other hand,

$$
\left(\overline{\mathcal{D}} A_{1} f\right)(z)=\int_{M^{\left(t_{0}\right)}} \frac{q(z)-q(w)}{(1-\langle z, w\rangle)^{2}} \overline{\left(\partial_{p_{z}} r\right)(z)}\left\langle z, p_{w}\right\rangle\left(\partial_{p_{w}} f\right)(w) d \mu(w)
$$

for $f \in \mathcal{P}$. Denote $\xi(z)=\left(\partial_{p_{z}} r\right)(z)$. Then the above translates to the operator identity

$$
\overline{\mathcal{D}} A_{1}=\sum_{j=1}^{n} M_{\chi_{M^{\left(t_{0}\right)}}} M_{\bar{\xi}} M_{\zeta_{j}}\left[M_{q}, \Lambda\right] M_{\bar{\psi}_{j}} D
$$

Applying Proposition 5.13, we conclude that $\overline{\mathcal{D}} A_{1}$ is in the trace class.
It is easy to see that

$$
\mathcal{I} A_{1}=\sum_{j=1}^{n} M_{\chi_{M^{\left(t_{0}\right)}}} M_{\zeta_{j}}\left[M_{\bar{r}},\left[M_{q}, \Lambda\right]\right] M_{\bar{\psi}_{j}} D .
$$

Applying Proposition 5.13 again, we see that $\mathcal{I} A_{1}$ is in the trace class. Lemma 6.9 now allows us to conclude that $\left[\hat{M}_{\bar{r}}, X_{1}\right] P=A_{1} \in \mathcal{C}_{1}$.

Let us now consider $A_{2}=\left[\hat{M}_{\bar{r}}, X_{2}\right] P$. For each $f \in \mathcal{P}$,

$$
\left(A_{2} f\right)(z)=\int_{M^{\left(t_{0}\right)}} \frac{(\bar{r}(z)-\bar{r}(w))\left(\partial_{p_{w}} q\right)(w)}{(1-\langle z, w\rangle)^{2}}\left\langle z, p_{w}\right\rangle f(w) d \mu(w)
$$

Thus

$$
\left(\mathcal{D} A_{2} f\right)(z)=\left(Y_{21} f\right)(z)+\left(Y_{22} f\right)(z)
$$

where

$$
\begin{aligned}
& \left(Y_{21} f\right)(z)=2 \int_{M^{\left(t_{0}\right)}} \frac{(\bar{r}(z)-\bar{r}(w))\left(\partial_{p_{w}} q\right)(w)}{(1-\langle z, w\rangle)^{3}}\left\langle p_{z}, w\right\rangle\left\langle z, p_{w}\right\rangle f(w) d \mu(w) \quad \text { and } \\
& \left(Y_{22} f\right)(z)=\int_{M^{\left(t_{0}\right)}} \frac{(\bar{r}(z)-\bar{r}(w))\left(\partial_{p_{w}} q\right)(w)}{(1-\langle z, w\rangle)^{2}}\left\langle p_{z}, p_{w}\right\rangle f(w) d \mu(w)
\end{aligned}
$$

It is easy to see that

$$
\begin{aligned}
& Y_{21}=2 \sum_{i=1}^{n} \sum_{j=1}^{n} M_{\chi_{M^{\left(t_{0}\right)}}} M_{\zeta_{j}} M_{\psi_{i}}\left[M_{\bar{r}}, Z\right] M_{\bar{\zeta}_{i}} M_{\bar{\psi}_{j}} M_{\eta} I \quad \text { and } \\
& Y_{22}=\sum_{j=1}^{n} M_{\chi_{M^{\left(t_{0}\right)}}} M_{\psi_{j}}\left[M_{\bar{r}}, \Lambda\right] M_{\bar{\psi}_{j}} M_{\eta} I .
\end{aligned}
$$

Thus it follows from Lemmas 5.3 and 6.3 that $Y_{21} \in \mathcal{C}_{1}$. Applying Proposition 5.13 again, we have $Y_{22} \in \mathcal{C}_{1}$. Hence $\mathcal{D} A_{2} \in \mathcal{C}_{1}$.

On the other hand,

$$
\left(\overline{\mathcal{D}} A_{2} f\right)(z)=\int_{M^{\left(t_{0}\right)}} \frac{1}{(1-\langle z, w\rangle)^{2}} \overline{\left(\partial_{p_{z}} r\right)(z)}\left(\partial_{p_{w}} q\right)(w)\left\langle z, p_{w}\right\rangle f(w) d \mu(w)
$$

$f \in \mathcal{P}$. Thus

$$
\overline{\mathcal{D}} A_{2}=\sum_{j=1}^{n} M_{\chi_{M^{\left(t_{0}\right)}}} M_{\bar{\xi}} M_{\zeta_{j}} \Lambda M_{\bar{\psi}_{j}} M_{\eta} I .
$$

Applying Lemmas 5.12 and 6.3 , we obtain the membership $\overline{\mathcal{D}} A_{2} \in \mathcal{C}_{1}$.
Finally, it is easy to see that

$$
\mathcal{I} A_{2}=\sum_{j=1}^{n} M_{\chi_{M^{\left(t_{0}\right)}}} M_{\zeta_{j}}\left[M_{\bar{r}}, \Lambda\right] M_{\bar{\psi}_{j}} M_{\eta} I
$$

Thus we have $\mathcal{I} A_{2} \in \mathcal{C}_{1}$ by Proposition 5.13. Having proved the memberships of $\mathcal{D} A_{2}$, $\overline{\mathcal{D}} A_{2}, \mathcal{I} A_{2}$ in $\mathcal{C}_{1}$, Lemma 6.9 allows us to conclude that $\left[\hat{M}_{\bar{r}}, X_{2}\right] P=A_{2} \in \mathcal{C}_{1}$. Since $X_{1}-X_{2}=\left[\hat{M}_{q}, \hat{T}_{1}\right]$, we have $\left[\hat{M}_{\bar{r}},\left[\hat{M}_{q}, \hat{T}_{1}\right]\right] P \in \mathcal{C}_{1}$ as promised.

## 7. Proof of Theorem 1.1

Recall from Section 4 that we write $M_{\zeta_{1}}, \ldots, M_{\zeta_{n}}$ for the restrictions of $\hat{M}_{\zeta_{1}}, \ldots$, $\hat{M}_{\zeta_{n}}$ to $\mathcal{P}$. Thus for each $1 \leq j \leq n, M_{\zeta_{j}}^{*}$ means the adjoint of $M_{\zeta_{j}}$ on $\mathcal{P}$.

Lemma 7.1. (a) For each $1 \leq j \leq n, M_{\zeta_{j}}^{*}=P \hat{M}_{\bar{\zeta}_{j}} P+I^{*} M_{\bar{\psi}_{j}} D$.
(b) For each $1 \leq j \leq n$, we have $M_{\zeta_{j}}^{*}-P \hat{M}_{\bar{\zeta}_{j}} P \in \mathcal{C}_{2}$.

Proof. Since $M_{\zeta_{j}}^{*}=P \hat{M}_{\zeta_{j}}^{*} P$ and $P \overline{\mathcal{D}}^{*}=(\overline{\mathcal{D}} P)^{*}=0$, (a) follows from Proposition 6.6. Lemma 6.3 tells us that $I \in \mathcal{C}_{2}$. Thus (b) follows from (a).
Proposition 7.2. For every pair of $i, j \in\{1, \ldots, n\}$, we have $\left[M_{\zeta_{i}}, M_{\zeta_{j}}^{*}-P \hat{M}_{\bar{\zeta}_{j}} P\right] \in \mathcal{C}_{1}$.
Proof. By Lemma 7.1, this is equivalent to the assertion that $\left[M_{\zeta_{i}}, I^{*} M_{\bar{\psi}_{j}} D\right] \in \mathcal{C}_{1}$ for all $i, j \in\{1, \ldots, n\}$. Note that $D M_{\zeta_{i}}=M_{\zeta_{i}} D+M_{\psi_{i}} I$. Therefore

$$
\left[M_{\zeta_{i}}, I^{*} M_{\bar{\psi}_{j}} D\right]=\left(M_{\zeta_{i}} I^{*}-I^{*} M_{\zeta_{i}}\right) M_{\bar{\psi}_{j}} D-I^{*} M_{\bar{\psi}_{j}} M_{\psi_{i}} I .
$$

It follows from Lemma 6.3 that $I^{*} M_{\bar{\psi}_{j}} M_{\psi_{i}} I \in \mathcal{C}_{1}$. Thus it suffices to show that $M_{\zeta_{i}} I^{*}-$ $I^{*} M_{\zeta_{i}}$ is in the trace class. Equivalently, it suffices to show that $I M_{\zeta_{i}}^{*}-M_{\bar{\zeta}_{i}} I$ is in the trace class. By Lemmas 7.1 and $6.3, I\left(M_{\zeta_{i}}^{*}-P \hat{M}_{\bar{\zeta}_{i}} P\right)$ is in the trace class. Hence our task is reduced to the proof that $I P \hat{M}_{\bar{\zeta}_{i}} P-M_{\bar{\zeta}_{i}} I$ is in the trace class.

Then note that $M_{\bar{\zeta}_{i}} I=M_{\bar{\zeta}_{i}} \mathcal{I} P=\mathcal{I} \hat{M}_{\bar{\zeta}_{i}} P$. Hence

$$
I P \hat{M}_{\bar{\zeta}_{i}} P-M_{\bar{\zeta}_{i}} I=-\mathcal{I}(1-P) \hat{M}_{\bar{\zeta}_{i}} P=-\mathcal{I}\left[\hat{M}_{\bar{\zeta}_{i}}, P\right] P
$$

To prove that this is in the trace class, we recall Proposition 4.6, which implies that $P=\hat{T}\left(\tilde{T}^{-1} \oplus 0\right)$. Thus

$$
\begin{equation*}
\mathcal{I}\left[\hat{M}_{\bar{\zeta}_{i}}, P\right] P=\mathcal{I}\left[\hat{M}_{\bar{\zeta}_{i}}, \hat{T}\right]\left(\tilde{T}^{-1} \oplus 0\right) P+\mathcal{I} \hat{T}\left[\hat{M}_{\bar{\zeta}_{i}}, \tilde{T}^{-1} \oplus 0\right] P \tag{7.1}
\end{equation*}
$$

Let us first consider the second term on the right-hand side. Since the range of $\hat{T}$ is contained in $\mathcal{P}$, it follows from Lemma 6.3 that $\mathcal{I} \hat{T}$ is a Hilbert-Schmidt operator. By Proposition 4.6, there is a $\xi \in C_{c}^{\infty}(\mathbf{R})$ such that $\tilde{T}^{-1} \oplus 0=\xi(\hat{T})$. Therefore it follows from Proposition 4.11 and smooth functional calculus that $\left[\hat{M}_{\bar{\zeta}_{i}}, \tilde{T}^{-1} \oplus 0\right] \in \mathcal{C}_{2}$. Thus the second term on the right-hand side of (7.1) is in the trace class.

The remaining task is to show that the first term on the right-hand side of (7.1) is in the trace class. Since $\hat{T}=\hat{T}_{1}+\hat{T}_{2}$, the proof of the proposition will be complete once we show that both $\mathcal{I}\left[\hat{M}_{\bar{\zeta}_{i}}, \hat{T}_{1}\right]$ and $\mathcal{I}\left[\hat{M}_{\bar{\zeta}_{i}}, \hat{T}_{2}\right]$ are in the trace class.

For this purpose, we define

$$
\begin{aligned}
& \left(T_{1}^{\circ} \varphi\right)(z)=\int_{M^{\left(t_{0}\right)}} \frac{\left\langle z, p_{w}\right\rangle}{(1-\langle z, w\rangle)^{2}} \varphi(w) d \mu(w) \quad \text { and } \\
& \left(T_{2}^{\circ} \varphi\right)(z)=\int_{M^{\left(t_{0}\right)}} \frac{1}{1-\langle z, w\rangle} \varphi(w) d \mu(w)
\end{aligned}
$$

$\varphi \in L^{2}\left(M^{\left(t_{0}\right)}, d \mu\right)$, which are operators on the space $L^{2}\left(M^{\left(t_{0}\right)}, d \mu\right)$. Note that $\mathcal{I} \hat{T}_{1}=T_{1}^{\circ} \mathcal{D}$. Since $\bar{\zeta}_{i}$ is conjugate analytic, we have $\mathcal{D} \hat{M}_{\bar{\zeta}_{i}}=M_{\bar{\zeta}_{i}} \mathcal{D}$. Also, $\mathcal{I} \hat{M}_{\bar{\zeta}_{i}}=M_{\bar{\zeta}_{i}} \mathcal{I}$. Thus

$$
\mathcal{I}\left[\hat{M}_{\bar{\zeta}_{i}}, \hat{T}_{1}\right]=\left[M_{\bar{\zeta}_{i}}, T_{1}^{\circ}\right] \mathcal{D}
$$

It is easy to see that

$$
\left[M_{\bar{\zeta}_{i}}, T_{1}^{\circ}\right]=\sum_{\ell=1}^{n} M_{\chi_{M^{\left(t_{0}\right)}}} M_{\zeta_{\ell}}\left[M_{\bar{\zeta}_{i}}, \Lambda\right] M_{\bar{\psi}_{\ell}} M_{\chi_{M^{\left(t_{0}\right)}}}
$$

Since $\left[M_{\bar{\zeta}_{i}}, \Lambda\right]=\left[\Lambda, M_{\zeta_{i}}\right]^{*}$, it follows from these identities and Proposition 5.13 that $\mathcal{I}\left[\hat{M}_{\bar{\zeta}_{i}}, \hat{T}_{1}\right]$ is in the trace class. It is also easy to see that

$$
\mathcal{I}\left[\hat{M}_{\bar{\zeta}_{i}}, \hat{T}_{2}\right]=\left[M_{\bar{\zeta}_{i}}, T_{2}^{\circ}\right] \mathcal{I}=M_{\chi_{M^{(t)}}}\left[M_{\bar{\zeta}_{i}}, \Lambda\right] \mathcal{I}-\sum_{\ell=1}^{n} M_{\chi_{M^{\left(t_{0}\right)}}} M_{\zeta_{\ell}}\left[M_{\bar{\zeta}_{i}}, \Lambda\right] M_{\bar{\zeta}_{\ell}} \mathcal{I}
$$

By Proposition 5.13, $\mathcal{I}\left[\hat{M}_{\widehat{\zeta}_{i}}, \hat{T}_{2}\right]$ is also in the trace class. This completes the proof.
After so much preparation, we are finally able to deal with essential normality.
Proposition 7.3. On the space $\mathcal{P}$, the commuting tuple $\left(M_{\zeta_{1}}, \ldots, M_{\zeta_{n}}\right)$ is 1-essentially normal. That is, for all $i, j \in\{1, \ldots, n\}$, we have $\left[M_{\zeta_{i}}, M_{\zeta_{j}}^{*}\right] \in \mathcal{C}_{1}$.

Proof. In view of Proposition 7.2, it suffices to show that $\left[M_{\zeta_{i}}, P \hat{M}_{\bar{\zeta}_{j}} P\right] \in \mathcal{C}_{1}$ for $i, j \in$ $\{1, \ldots, n\}$. Since $\hat{M}_{\zeta_{i}}$ commutes with $\hat{M}_{\bar{\zeta}_{j}}$, we have

$$
\begin{aligned}
{\left[M_{\zeta_{i}}, P \hat{M}_{\bar{\zeta}_{j}} P\right] } & =\left[P \hat{M}_{\zeta_{i}} P, P \hat{M}_{\bar{\zeta}_{j}} P\right]=P \hat{M}_{\bar{\zeta}_{j}}(1-P) \hat{M}_{\zeta_{i}} P-P \hat{M}_{\zeta_{i}}(1-P) \hat{M}_{\bar{\zeta}_{j}} P \\
& =\left[P, \hat{M}_{\bar{\zeta}_{j}}\right](1-P)\left[\hat{M}_{\zeta_{i}}, P\right]-\left[P, \hat{M}_{\zeta_{i}}\right](1-P)\left[\hat{M}_{\bar{\zeta}_{j}}, P\right] .
\end{aligned}
$$

By Corollary 6.11, this is in the trace class.
Lemma 7.4. We have $\left[M_{\zeta_{i}}, \tilde{T}\right] \in \mathcal{C}_{2}$ for every $1 \leq i \leq n$.
Proof. Obviously, we have $\left[M_{\zeta_{i}}, \tilde{T}\right]=P\left[\hat{M}_{\zeta_{i}}, \hat{T}\right] P$. Thus the membership $\left[M_{\zeta_{i}}, \tilde{T}\right] \in \mathcal{C}_{2}$ is a consequence of Proposition 4.11.
Proposition 7.5. For all $i, j \in\{1, \ldots, n\}$, we have $\left[M_{\zeta_{i}},\left[M_{\zeta_{j}}^{*}, \tilde{T}\right]\right] \in \mathcal{C}_{1}$.
Proof. It follows from Proposition 7.2 and Lemmas 7.1 and 7.4 that $\left[M_{\zeta_{i}},\left[M_{\zeta_{j}}^{*}-P \hat{M}_{\bar{\zeta}_{j}} P, \tilde{T}\right]\right]$ $\in \mathcal{C}_{1}$ for all $i, j \in\{1, \ldots, n\}$. Therefore it suffices to show that $\left[M_{\zeta_{i}},\left[P \hat{M}_{\bar{\zeta}_{j}} P, \tilde{T}\right]\right] \in \mathcal{C}_{1}$, $i, j \in\{1, \ldots, n\}$. Again, let us write $A \sim_{1} B$ when $A-B \in \mathcal{C}_{1}$. We have

$$
\left[M_{\zeta_{i}},\left[P \hat{M}_{\bar{\zeta}_{j}} P, \tilde{T}\right]\right]=P\left[\hat{M}_{\zeta_{i}}, P\left[\hat{M}_{\bar{\zeta}_{j}}, \hat{T}\right] P\right] P \sim_{1} P\left[\hat{M}_{\zeta_{i}},\left[\hat{M}_{\bar{\zeta}_{j}}, \hat{T}\right]\right] P=P\left[\hat{M}_{\bar{\zeta}_{j}},\left[\hat{M}_{\zeta_{i}}, \hat{T}\right]\right] P
$$

where the $\sim_{1}$ follows from Proposition 4.11 and Corollary 6.11. Now an application of Proposition 4.12 completes the proof.
Corollary 7.6. For all $i, j \in\{1, \ldots, n\}$, we have $\left[M_{\zeta_{i}},\left[M_{\zeta_{j}}^{*}, \tilde{T}^{1 / 2}\right]\right] \in \mathcal{C}_{1}$.
Proof. By Proposition 4.6, there are $0<c<C<\infty$ such that the spectrum of $\tilde{T}$ is contained in the interval $[c, C]$. Consider $H_{+}=\{\lambda \in \mathbf{C}: \operatorname{Re}(\lambda)>0\}$, the right half-plane. Let $\gamma$ be a simple Jordan curve in $H_{+} \backslash[c, C]$ whose winding number about every $x \in[c, C]$ is 1 . Taking advantage of the fact that the square-root function $\lambda^{1 / 2}$ is analytic on $H_{+}$, from the Riesz functional calculus we obtain

$$
\begin{equation*}
\tilde{T}^{1 / 2}=\frac{1}{2 \pi i} \int_{\gamma} \lambda^{1 / 2}(\lambda-\tilde{T})^{-1} d \lambda \tag{7.2}
\end{equation*}
$$

Therefore

$$
\left[M_{\zeta_{i}},\left[M_{\zeta_{j}}^{*}, \tilde{T}^{1 / 2}\right]\right]=\frac{1}{2 \pi i} \int_{\gamma} \lambda^{1 / 2}\{A(\lambda)+B(\lambda)+C(\lambda)\} d \lambda,
$$

where

$$
\begin{aligned}
& A(\lambda)=(\lambda-\tilde{T})^{-1}\left[M_{\zeta_{i}}, \tilde{T}\right](\lambda-\tilde{T})^{-1}\left[M_{\zeta_{j}}^{*}, \tilde{T}\right](\lambda-\tilde{T})^{-1} \\
& B(\lambda)=(\lambda-\tilde{T})^{-1}\left[M_{\zeta_{i}},\left[M_{\zeta_{j}}^{*}, \tilde{T}\right]\right](\lambda-\tilde{T})^{-1} \quad \text { and } \\
& C(\lambda)=(\lambda-\tilde{T})^{-1}\left[M_{\zeta_{j}}^{*}, \tilde{T}\right](\lambda-\tilde{T})^{-1}\left[M_{\zeta_{i}}, \tilde{T}\right](\lambda-\tilde{T})^{-1}
\end{aligned}
$$

Applying Proposition 7.5 to $B(\lambda)$ and Lemma 7.4 to $A(\lambda)$ and $C(\lambda)$, we obtain the membership $\left[M_{\zeta_{i}},\left[M_{\zeta_{j}}^{*}, \tilde{T}^{1 / 2}\right]\right] \in \mathcal{C}_{1}$.
Corollary 7.7. We have $\left[M_{\zeta_{i}}, \tilde{T}^{1 / 2}\right] \in \mathcal{C}_{2}$ for every $1 \leq i \leq n$.
Proof. This follows obviously from Lemma 7.4 and identity (7.2).
Finally, we can accomplish the main goal of the paper:
Proof of Theorem 1.1. Recall from Section 4 that the operator $J: \mathcal{Q} \rightarrow \mathcal{P}$ is invertible. Therefore the operator $J^{*}: \mathcal{P} \rightarrow \mathcal{Q}$ is also invertible. By the standard polarization,

$$
J^{*}=U\left|J^{*}\right|
$$

where $U: \mathcal{P} \rightarrow \mathcal{Q}$ is a unitary operator and $\left|J^{*}\right|=\left(J J^{*}\right)^{1 / 2}=\tilde{T}^{1 / 2}$ (see Proposition 4.6). Recall from Corollary 4.10 that $J Q_{\zeta_{j}}=M_{\zeta_{j}} J$ for every $j \in\{1, \ldots, n\}$. Therefore

$$
Q_{\zeta_{j}}=J^{-1} M_{\zeta_{j}} J=U \tilde{T}^{-1 / 2} M_{\zeta_{j}} \tilde{T}^{1 / 2} U^{*}=U M_{\zeta_{j}} U^{*}+K_{j},
$$

$j \in\{1, \ldots, n\}$, where

$$
K_{j}=U \tilde{T}^{-1 / 2}\left[M_{\zeta_{j}}, \tilde{T}^{1 / 2}\right] U^{*}
$$

Thus for any $i, j \in\{1, \ldots, n\}$, we have

$$
\left[Q_{\zeta_{i}}, Q_{\zeta_{j}}^{*}\right]=U\left[M_{\zeta_{i}}, M_{\zeta_{j}}^{*}\right] U^{*}+\left[U M_{\zeta_{i}} U^{*}, K_{j}^{*}\right]+\left[K_{i}, U M_{\zeta_{j}}^{*} U^{*}\right]+\left[K_{i}, K_{j}^{*}\right] .
$$

By Proposition 7.3 and Corollary 7.7, the first term and the last term on the right-hand side are in the trace class. What remains is to show that the two middle terms on the right-hand side are also in the trace class.

We have

$$
\begin{aligned}
& {\left[U M_{\zeta_{i}} U^{*},\right.} \\
& \left.\quad K_{j}^{*}\right]=U\left[M_{\zeta_{i}},\left\{\tilde{T}^{-1 / 2}\left[M_{\zeta_{j}}, \tilde{T}^{1 / 2}\right]\right\}^{*}\right] U^{*}=U\left[M_{\zeta_{i}},\left[\tilde{T}^{1 / 2}, M_{\zeta_{j}}^{*}\right] \tilde{T}^{-1 / 2}\right] U^{*} \\
& \quad=U\left[M_{\zeta_{i}},\left[\tilde{T}^{1 / 2}, M_{\zeta_{j}}^{*}\right]\right] \tilde{T}^{-1 / 2} U^{*}+U\left[\tilde{T}^{1 / 2}, M_{\zeta_{j}}^{*}\right] \tilde{T}^{-1 / 2}\left[\tilde{T}^{1 / 2}, M_{\zeta_{i}}\right] \tilde{T}^{-1 / 2} U^{*} \\
& \quad=A+B .
\end{aligned}
$$

By Corollary 7.6, we have $A \in \mathcal{C}_{1}$. By Corollary 7.7, we have $B \in \mathcal{C}_{1}$. Hence [ $U M_{\zeta_{i}} U^{*}, K_{j}^{*}$ ] $\in \mathcal{C}_{1}$. Since $\left[K_{i}, U M_{\zeta_{j}}^{*} U^{*}\right]=\left[U M_{\zeta_{j}} U^{*}, K_{i}^{*}\right]^{*}$, we also have $\left[K_{i}, U M_{\zeta_{j}}^{*} U^{*}\right] \in \mathcal{C}_{1}$. This completes the proof of Theorem 1.1.

## Appendix 1

The purpose of this appendix is to give a proof for Lemma 2.13. With minor changes of details, this is essentially the same as the proof of [29, Lemma 3.1].

Let $-1<\tau<\infty$ be given. Then for any $f \in C_{c}^{1}[0, \infty)$, we have

$$
\begin{equation*}
\int_{0}^{\infty}|f(x)|^{2} x^{\tau} d x \leq\left(\frac{2}{\tau+1}\right)^{2} \int_{0}^{\infty}\left|x f^{\prime}(x)\right|^{2} x^{\tau} d x \tag{A1.1}
\end{equation*}
$$

See inequality (3.1) in [29].
Let $0<b<t \leq 1$ be such that $M^{(b)} \subset \mathcal{M}$. Consider the function $\rho(w)=1-|w|^{2}$. Since $\mathcal{M}$ intersects $S=\left\{\zeta \in \mathbf{C}^{n}:|\zeta|=1\right\}$ transversely, $\nabla_{\mathcal{M}} \rho$ does not vanish near $\mathcal{M} \cap S$. Thus we can use $\rho$ as one of the real coordinates on $\mathcal{M}$ near $S$. More precisely, if $\zeta \in \mathcal{M} \cap S$, then $\zeta$ has an open neighborhood $U_{\zeta}$ in $\mathcal{M}$ that has the following properties:
(1) $U_{\zeta}=G((-c, c) \times V)$, where $0<c<b, V$ is a bounded open set in $\mathbf{R}^{2 d-1}$ and
$G:(-c, c) \times V \rightarrow \mathbf{C}^{n}$ is a one-to-one $C^{\infty}$ map.
(2) there are $0<\delta<C<\infty$ such that $D G$, the derivative of $G$, satisfies the matrix inequality $\delta \leq(D G)^{*}(x, y)(D G)(x, y) \leq C$ for all $x \in(-c, c)$ and $y \in V$.
(3) If $w=G(x, y)$ for some $x \in(-c, c)$ and $y \in V$, then $x=1-|w|^{2}$. Equivalently,
for each $w \in U_{\zeta}$, there is a unique $y_{w} \in V$ such that $w=G\left(1-|w|^{2}, y_{w}\right)$.
Obviously, (3) implies $U_{\zeta} \cap M=G((0, c) \times V) \subset M^{(b)}$.
Once we have this $c$, by the standard technique of using a smooth cutoff function, we can truncate inequality (A1.1) to the interval $[0, c]$. That is, there are $C_{1}$ and $C_{2}$ such that

$$
\begin{equation*}
\int_{0}^{c}|h(x)|^{2} x^{\tau} d x \leq C_{1} \int_{0}^{c}\left|x h^{\prime}(x)\right|^{2} x^{\tau} d x+C_{2} \int_{c / 2}^{c}|h(x)|^{2} x^{\tau} d x \tag{A1.2}
\end{equation*}
$$

for every $h \in C^{1}[0, c]$. Let $f$ be any $C^{1}$ function on an open set containing the closure of $M^{(t)}$. By the definition of $d \nu_{\tau}$ and property (3) above,

$$
\begin{aligned}
& \int_{U_{\zeta} \cap M}|f(w)|^{2} d \nu_{\tau}(w)=\int_{U_{\zeta} \cap M}|f(w)|^{2}\left(1-|w|^{2}\right)^{\tau} d v_{M}(w) \\
& \quad=\int_{V} \int_{0}^{c}|f(G(x, y))|^{2} x^{\tau} J(x, y) d x d y \leq C_{3} \int_{V} \int_{0}^{c}|f(G(x, y))|^{2} x^{\tau} d x d y \\
& \quad \leq C_{4} \int_{V} \int_{0}^{c}\left|x \frac{d}{d x} f(G(x, y))\right|^{2} x^{\tau} d x d y+C_{5} \int_{V} \int_{c / 2}^{c}|f(G(x, y))|^{2} x^{\tau} d x d y
\end{aligned}
$$

where the second $\leq$ is an application of (A1.2). By the chain rule for differentiation,

$$
\frac{d}{d x} f(G(x, y))=\langle(\nabla f)(G(x, y)), \tau(x, y)\rangle
$$

where $\tau(x, y)$ is a (real) tangent vector to $\mathcal{M}$ at the point $G(x, y)$. Moreover, (2) implies the bound $|\tau(x, y)| \leq C^{1 / 2}$. Hence $|d f(G(x, y)) / d x| \leq C^{1 / 2}\left|\left(\nabla_{\mathcal{M}} f\right)(G(x, y))\right|$. Thus

$$
\begin{aligned}
& \int_{V} \int_{0}^{c} \mid\left.x \frac{d}{d x} f(G(x, y))\right|^{2} x^{\tau} d x d y \leq C \int_{V} \int_{0}^{c}\left|\left(\nabla_{\mathcal{M}} f\right)(G(x, y))\right|^{2} x^{2+\tau} d x d y \\
& \leq C_{6} \int_{V} \int_{0}^{c}\left|\left(\nabla_{\mathcal{M}} f\right)(G(x, y))\right|^{2} x^{2+\tau} J(x, y) d x d y \\
& \quad=C_{6} \int_{U_{\zeta} \cap M}\left|\left(\nabla_{\mathcal{M}} f\right)(w)\right|^{2}\left(1-|w|^{2}\right)^{2+\tau} d v_{M}(w) \\
& \quad=C_{6} \int_{U_{\zeta} \cap M}\left|\left(\nabla_{\mathcal{M}} f\right)(w)\right|^{2}\left(1-|w|^{2}\right)^{2} d \nu_{\tau}(w)
\end{aligned}
$$

where the third step uses property (3). Combining this with (A1.3), we find that

$$
\begin{align*}
\int_{U_{\zeta} \cap M}|f(w)|^{2} d \nu_{\tau}(w) & \leq C_{7} \int_{U_{\zeta} \cap M}\left|\left(\nabla_{\mathcal{M}} f\right)(w)\right|^{2}\left(1-|w|^{2}\right)^{2} d \nu_{\tau}(w) \\
& +C_{8} \int_{N^{(c / 2)} \cap M^{(t)}}|f(w)|^{2} d \nu_{\tau}(w) \tag{A1.4}
\end{align*}
$$

Since $\mathcal{M} \cap S$ is compact, there are $\zeta_{1}, \ldots, \zeta_{k} \in \mathcal{M} \cap S$ such that the corresponding open sets $U_{\zeta_{1}}, \ldots, U_{\zeta_{k}}$ have the property $U_{\zeta_{1}} \cup \cdots \cup U_{\zeta_{k}} \supset \mathcal{M} \cap S=\left\{w \in \mathcal{M}: 1-|w|^{2}=0\right\}$. Thus $U_{\zeta_{1}} \cup \cdots \cup U_{\zeta_{k}} \supset M^{(\sigma)}$ for some $0<\sigma<1$. Combining this containment with (A1.4), Lemma 2.13 is proved.

## Appendix 2

The purpose of this appendix is to give a proof of Proposition 5.6. One can characterize the material in this appendix as an easier version of parts of Sections 3 and 4.

Recall that $\mathcal{H}^{(2-n)}$ is the Hilbert space of analytic functions on $\mathbf{B}$ which has (5.2) as its reproducing kernel. We begin with an easier version of Lemma 3.9:

Lemma A2.1. Given any $\delta>0$, there exist $0<\tau<t_{0}$ and $0<C<\infty$ such that

$$
\int_{M^{(t)}}|f(z)|^{2} d \mu(z) \leq C \int_{M^{(t)}}|(Z f)(z)|^{2} d \mu(z)+\delta \int_{M^{\left(t_{0}\right)}}|f(z)|^{2} d \mu(z)
$$

for all $0<t \leq \tau$ and $f \in \mathcal{H}^{(2-n)}$.
Proof. As in the proof of Lemma 3.9, we begin with a large $1 \leq r<\infty$, whose exact value will be determined later. For such an $r$, let $0<\tau_{3} \leq \tau_{2} \leq \tau_{1} \leq t_{0}$ be the same as in the proof of Lemma 3.9. Thus if $0<t \leq \tau_{2}$, then

$$
\begin{equation*}
I_{z}\left(D(z, 2 r) \cap T_{z}^{\bmod }\right) \supset D(z, r) \cap M \quad \text { for every } \quad z \in M^{(t)} \tag{A2.1}
\end{equation*}
$$

And if $0<t \leq \tau_{3}$, then

$$
\begin{equation*}
I_{z}\left(D(z, 2 r) \cap T_{z}^{\bmod }\right) \subset M^{\left(t_{0}\right)} \quad \text { for every } \quad z \in M^{(t)} \tag{A2.2}
\end{equation*}
$$

Write $U(z)=I_{z}\left(D(z, 2 r) \cap T_{z}^{\bmod }\right)$ as before, $z \in M^{\left(\tau_{3}\right)}$. Let $f \in \mathcal{H}^{(2-n)}$ be given. Then

$$
(Z f)(z)=A(z)+B(z)
$$

where

$$
\begin{aligned}
& A(z)=\int_{U(z)} f(w) \frac{1-|w|^{2}}{(1-\langle z, w\rangle)^{3}} d v_{M}(w) \quad \text { and } \\
& B(z)=\int_{M^{\left(t_{0}\right)} \backslash U(z)} f(w) \frac{1-|w|^{2}}{(1-\langle z, w\rangle)^{3}} d v_{M}(w)
\end{aligned}
$$

$z \in M^{\left(\tau_{3}\right)}$. Since $P_{z} U(z)=D(z, 2 r) \cap T_{z}^{\bmod }$, by (2.10) we have

$$
A(z)=\int_{D(z, 2 r) \cap T_{z}^{\bmod }} f\left(I_{z}(x)\right) \frac{1-\left|I_{z}(x)\right|^{2}}{\left(1-\left\langle z, I_{z}(x)\right\rangle\right)^{3}} J_{z}(x) d v_{1}(x)
$$

Recall from (2.6) that $\left\langle z, I_{z}(x)\right\rangle=\langle z, x\rangle$. Writing

$$
F(z, x)=1-\frac{1-|x|^{2}}{1-\left|I_{z}(x)\right|^{2}} \cdot \frac{J_{z}(z)}{J_{z}(x)}
$$

we have $A(z)=A_{1}(z)+A_{2}(z)$, where

$$
\begin{aligned}
& A_{1}(z)=J_{z}(z) \int_{D(z, 2 r) \cap T_{z}^{\bmod }} f\left(I_{z}(x)\right) \frac{1-|x|^{2}}{(1-\langle z, x\rangle)^{3}} d v_{1}(x) \quad \text { and } \\
& A_{2}(z)=\int_{D(z, 2 r) \cap T_{z}^{\bmod }} f\left(I_{z}(x)\right) \frac{1-\left|I_{z}(x)\right|^{2}}{\left(1-\left\langle z, I_{z}(x)\right\rangle\right)^{3}} F(z, x) J_{z}(x) d v_{1}(x)
\end{aligned}
$$

Let us analyze $A_{1}(z)$, and $A_{2}(z)$ and $B(z)$.
Being a local inverse of $P_{z}$, the map $I_{z}$ is analytic. Thus the function $x \mapsto f\left(I_{z}(x)\right)$ is analytic on $D(z, 3 r) \cap T_{z}^{\bmod }$. Therefore it follows from Lemma 2.10 that

$$
A_{1}(z)=\Phi(2 r) J_{z}(z) f\left(I_{z}(z)\right)=\Phi(2 r) J_{z}(z) f(z)
$$

Recalling (2.9) and (2.7), we see that there is a $0<C_{1}<\infty$ such that

$$
\begin{equation*}
|f(z)| \leq C_{1}\left|A_{1}(z)\right| \quad \text { for } \quad z \in M^{(t)}, \quad 0<t \leq \tau_{3} \tag{A2.3}
\end{equation*}
$$

Denote

$$
\epsilon(r, t)=\sup _{z \in M^{(t)}}\left\{\sup _{x \in D(z, 2 r) \cap T_{z}^{\bmod }}|F(z, x)|\right\} .
$$

As we explained in the proof of Lemma 3.9,

$$
\begin{equation*}
\lim _{t \downarrow 0} \epsilon(r, t)=0 \tag{A2.4}
\end{equation*}
$$

for every $1 \leq r<\infty$. Applying (2.10) again, we have

$$
\begin{aligned}
\left|A_{2}(z)\right| & \leq \epsilon(r, t) \int_{D(z, 2 r) \cap T_{z}^{\bmod }}\left|f\left(I_{z}(x)\right)\right| \frac{1-\left|I_{z}(x)\right|^{2}}{\left|1-\left\langle z, I_{z}(x)\right\rangle\right|^{3}} J_{z}(x) d v(x) \\
& \leq \epsilon(r, t) \int_{M^{\left(t_{0}\right)}}|f(w)| \frac{1}{|1-\langle z, w\rangle|^{3}} d \mu(w) .
\end{aligned}
$$

Thus it follows from Lemma 3.8 that

$$
\begin{equation*}
\int_{M^{(t)}}\left|A_{2}(z)\right|^{2} d \mu(z) \leq\{\epsilon(r, t)\}^{2}\|B\|^{2} \int_{M^{\left(t_{0}\right)}}|f(w)|^{2} d \mu(w) . \tag{A2.5}
\end{equation*}
$$

Finally, from (A2.1) we obtain

$$
|B(z)| \leq \int_{M^{\left(t_{0}\right)} \backslash D(z, r)}|f(w)| \frac{1}{|1-\langle z, w\rangle|^{3}} d \mu(w)
$$

for $z \in M^{(t)}, 0<t \leq \tau_{3}$. Using the operator $B_{r}$ in Lemma 3.8, we have

$$
\begin{equation*}
\int_{M^{(t)}}|B(z)|^{2} d \mu(z) \leq\left\|B_{r}\right\|^{2} \int_{M^{\left(t_{0}\right)}}|f(w)|^{2} d \mu(w) \tag{A2.6}
\end{equation*}
$$

Retracing the above steps, we have

$$
A_{1}(z)=(Z f)(z)-\left(A_{2}(z)+B(z)\right)
$$

Thus it follows from (A2.3), (A2.5) and (A2.6) that

$$
\begin{align*}
\int_{M^{(t)}}|f(z)|^{2} d \mu(z) & \leq 3 C_{1}^{2} \int_{M^{(t)}}|(Z f)(z)|^{2} d \mu(z) \\
& +3 C_{1}^{2}\left(\{\epsilon(r, t)\}^{2}\|B\|^{2}+\left\|B_{r}\right\|^{2}\right) \int_{M^{\left(t_{0}\right)}}|f(z)|^{2} d \mu(z) \tag{A2.7}
\end{align*}
$$

Let any $\delta>0$ be given. By Lemma 3.8, we can first pick an $r \in[1, \infty)$ such that $3 C_{1}^{2}\left\|B_{r}\right\|^{2} \leq \delta / 2$. With $r$ so fixed, by (A2.4), we can pick a $0<\tau \leq \tau_{3}$ such that

$$
3 C_{1}^{2}\{\epsilon(r, t)\}^{2}\|B\|^{2} \leq \delta / 2
$$

for every $0<t \leq \tau$. Substitution these bounds in (A2.7), the lemma is proved.
On the space $\mathcal{H}^{(2-n)}$, we now define the operator

$$
T^{\dagger}=\int_{M} K_{w}^{(2-n)} \otimes K_{w}^{(2-n)} d \mu(w)
$$

By (5.2) and the reproducing property of $K_{w}^{(2-n)}$, we have

$$
\begin{equation*}
\left(T^{\dagger} f\right)(z)=(Z f)(z) \quad \text { for } \quad z \in M \tag{A2.8}
\end{equation*}
$$

$f \in \mathcal{H}^{(2-n)}$. Write $\langle\cdot, \cdot\rangle_{2-n}$ for the inner product on $\mathcal{H}^{(2-n)}$. Then

$$
\begin{equation*}
\left\langle T^{\dagger} f, f\right\rangle_{2-n}=\int_{M}|f(w)|^{2} d \mu(w) \tag{A2.9}
\end{equation*}
$$

$f \in \mathcal{H}^{(2-n)}$. Thus it follows from Lemma 5.14 that

$$
\left\langle T^{\dagger} f, f\right\rangle_{2-n} \leq C\|f\|_{2-n}^{2} \quad \text { for every } \quad f \in \mathcal{H}^{(2-n)}
$$

In other words, $T^{\dagger}$ is a bounded operator on $\mathcal{H}^{(2-n)}$. Obviously, $T^{\dagger} \geq 0$.
Lemma A2.2. For each $y \in \mathcal{M} \cap M=\mathcal{M} \cap \mathbf{B}$, there is an open neighborhood $N_{y}$ of $y$ in $\mathcal{M} \cap M$ which has the following property. Let $\left\{f_{k}\right\}$ be a sequence in $\mathcal{H}^{(2-n)}$. If the sequence $\left\{\left(T^{\dagger}\right)^{1 / 2} f_{k}\right\}$ weakly converges to 0 , then

$$
\lim _{k \rightarrow \infty} \sup \left\{\left|f_{k}(w)\right|: w \in N_{y}\right\}=0
$$

Proof. This is very similar to the proof of Lemma 3.7. Indeed we again use the biholomorphic map $\rho_{y}: B_{1}(0,2) \rightarrow V_{y}$ introduced in Section 3, $y \in \mathcal{M}$. Recall that $\rho_{y}(0)=y$. For each $y \in \mathcal{M} \cap M$, there are $\epsilon=\epsilon(y)>0$ and $r=r(y) \in(0,2)$ such that $\rho_{y}\left(B_{1}(0, r)\right) \subset\left\{w \in \mathcal{M} \cap M: 1-|w|^{2}>\epsilon\right\}$. We will show that the lemma holds for the open set $N_{y}=\rho_{y}\left(B_{1}(0, r / 2)\right)$.

Again, consider the Bergman space $L_{a}^{2}\left(B_{1}(0, r), d v_{1}\right)$. This time, we define

$$
(G f)(\xi)=f\left(\rho_{y}(\xi)\right), \quad \xi \in B_{1}(0, r)
$$

$f \in \mathcal{H}^{(2-n)}$. By the condition $r<2$, we have

$$
\begin{array}{rl}
\int_{B_{1}(0, r)}|(G f)(\xi)|^{2} & d v_{1}(\xi) \leq C_{1} \int_{B_{1}(0, r)}\left|f\left(\rho_{y}(\xi)\right)\right|^{2}\left|\rho_{y}^{\prime}(\xi)\right|^{2} d v_{1}(\xi) \\
& =C_{1} \int_{\rho_{y}\left(B_{1}(0, r)\right)}|f(w)|^{2} d v_{M}(w) \\
& \leq C_{1} \epsilon^{-1} \int_{\rho_{y}\left(B_{1}(0, r)\right)}|f(w)|^{2}\left(1-|w|^{2}\right) d v_{M}(w) \\
& =C_{1} \epsilon^{-1} \int_{\rho_{y}\left(B_{1}(0, r)\right)}|f(w)|^{2} d \mu(w) \\
& \leq C_{1} \epsilon^{-1}\left\langle T^{\dagger} f, f\right\rangle_{2-n}=C_{1} \epsilon^{-1}\left\|\left(T^{\dagger}\right)^{1 / 2} f\right\|_{2-n}^{2}
\end{array}
$$

Thus $G=W\left(T^{\dagger}\right)^{1 / 2}$, where $W: \mathcal{H}^{(2-n)} \rightarrow L_{a}^{2}\left(B_{1}(0, r), d v_{1}\right)$ is a bounded operator.
Now let $\left\{f_{k}\right\}$ be any sequence in $\mathcal{H}^{(2-n)}$ such that $\left\{\left(T^{\dagger}\right)^{1 / 2} f_{k}\right\}$ weakly converges to 0 . Since $G=W\left(T^{\dagger}\right)^{1 / 2}$, the sequence $\left\{G f_{k}\right\}$ weakly converges to 0 in $L_{a}^{2}\left(B_{1}(0, r), d v_{1}\right)$. Using the reproducing kernel for the Bergman space, we have

$$
\lim _{k \rightarrow \infty} \sup \left\{\left|f_{k}\left(\rho_{y}(\xi)\right)\right|: \xi \in B_{1}(0, r / 2)\right\}=0
$$

Since $N_{y}=\rho_{y}\left(B_{1}(0, r / 2)\right)$, the proof is complete.
Proposition A2.3. There is a $\gamma>0$ such that the spectrum of $T^{\dagger}$ does not intersect the interval $(0, \gamma)$.

Proof. By Lemma A2.1, there are $0<t<t_{0}$ and $0<C<\infty$ such that

$$
\int_{M^{(t)}}|f(z)|^{2} d \mu(z) \leq C \int_{M^{(t)}}|(Z f)(z)|^{2} d \mu(z)+\frac{1}{2} \int_{M^{\left(t_{0}\right)}}|f(z)|^{2} d \mu(z)
$$

for every $f \in \mathcal{H}^{(2-n)}$. After the obvious cancellation, we have

$$
\frac{1}{2} \int_{M^{(t)}}|f(z)|^{2} d \mu(z) \leq C \int_{M^{(t)}}|(Z f)(z)|^{2} d \mu(z)+\frac{1}{2} \int_{M^{\left(t_{0}\right)} \backslash M^{(t)}}|f(z)|^{2} d \mu(z)
$$

Combining this with (A2.8) and (A2.9), we find that

$$
\begin{equation*}
\frac{1}{2} \int_{M^{(t)}}|f(z)|^{2} d \mu(z) \leq C\left\|\left(T^{\dagger}\right)^{3 / 2} f\right\|_{2-n}^{2}+\frac{1}{2} \int_{M^{\left(t_{0}\right)} \backslash M^{(t)}}|f(z)|^{2} d \mu(z) \tag{A2.10}
\end{equation*}
$$

for every $f \in \mathcal{H}^{(2-n)}$.
Let $d E$ be the spectral measure for the positive operator $T^{\dagger}$. Suppose that $E(0, \gamma) \neq 0$ for every $\gamma>0$. We will complete the proof by showing that this leads to a contradiction. For each $k \in \mathbf{N}$, since $E(0,1 / k) \neq 0$, we pick an $f_{k} \in E(0,1 / k) \mathcal{H}^{(2-n)}$ such that $\left\langle T^{\dagger} f_{k}, f_{k}\right\rangle_{2-n}=1$. By (A2.9). this means

$$
\begin{equation*}
\int_{M}\left|f_{k}(w)\right|^{2} d \mu(w)=1 \tag{A2.11}
\end{equation*}
$$

for every $k$. Obviously, the sequence $\left\{\left(T^{\dagger}\right)^{1 / 2} f_{k}\right\}$ weakly converges to 0 in $\mathcal{H}^{(2-n)}$.
Let $H$ be the closure of $M^{\left(t_{0}\right)} \backslash M^{(t)}$. Then $H$ is a compact subset of $\mathcal{M} \cap M$. By Lemma A2.2 and a usual covering argument, the weak convergence to 0 of the sequence $\left\{\left(T^{\dagger}\right)^{1 / 2} f_{k}\right\}$ implies

$$
\lim _{k \rightarrow \infty} \sup \left\{\left|f_{k}(z)\right|: z \in H\right\}=0
$$

By the maximum modulus principle, this implies

$$
\lim _{k \rightarrow \infty} \sup \left\{\left|f_{k}(z)\right|: z \in M \backslash M^{(t)}\right\}=0
$$

Therefore

$$
\begin{equation*}
\lim _{k \rightarrow \infty} \int_{M \backslash M^{(t)}}\left|f_{k}(z)\right|^{2} d \mu(z)=0 \tag{A2.12}
\end{equation*}
$$

Combining this with (A2.11), we find that

$$
\begin{equation*}
\lim _{k \rightarrow \infty} \int_{M^{(t)}}\left|f_{k}(z)\right|^{2} d \mu(z)=1 \tag{A2.13}
\end{equation*}
$$

Since $f_{k} \in E(0,1 / k) \mathcal{H}^{(2-n)}$, we have

$$
\begin{equation*}
\left\|\left(T^{\dagger}\right)^{3 / 2} f\right\|_{2-n}^{2} \leq k^{-2}\left\langle T^{\dagger} f_{k}, f_{k}\right\rangle=k^{-2} \tag{A2.14}
\end{equation*}
$$

$k \in \mathbf{N}$. On the other hand, specializing (A2.10) to each $f_{k}$, we see that

$$
\begin{equation*}
\frac{1}{2} \int_{M^{(t)}}\left|f_{k}(z)\right|^{2} d \mu(z) \leq C\left\|\left(T^{\dagger}\right)^{3 / 2} f_{k}\right\|_{2-n}^{2}+\frac{1}{2} \int_{M^{\left(t_{0}\right)} \backslash M^{(t)}}\left|f_{k}(z)\right|^{2} d \mu(z) \tag{A2.15}
\end{equation*}
$$

$k \in \mathbf{N}$. Obviously, the combination of (A2.15) with (A2.12), (A2.13) and (A2.14) leads to the contradiction $1 / 2 \leq 0$. This completes the proof.

Define

$$
\mathcal{R}^{(2-n)}=\left\{f \in \mathcal{H}^{(2-n)}: f=0 \quad \text { on } M\right\} .
$$

Then $\mathcal{R}^{(2-n)}$ is a submodule of the Hilbert module $\mathcal{H}^{(2-n)}$, just as $\mathcal{R}$ is a submodule of the Drury-Arveson module $H_{n}^{2}$.

Proof of Proposition 5.6. Define

$$
\mathcal{Q}^{(2-n)}=\mathcal{H}^{(2-n)} \ominus \mathcal{R}^{(2-n)},
$$

which is the quotient module of $\mathcal{H}^{(2-n)}$ corresponding to the submodule $\mathcal{R}^{(2-n)}$. Let $J^{(2-n)}$ denote the operator that takes each $f \in \mathcal{H}^{(2-n)}$ to the same $f$ in $L^{2}(M, d \mu)$. That is, $J^{(2-n)} f$ is the restriction of $f \in \mathcal{H}^{(2-n)}$ to the subset $M$ of $\mathbf{B}$. Then (A2.9) translates to the operator identity

$$
\begin{equation*}
\left(J^{(2-n)}\right)^{*} J^{(2-n)}=T^{\dagger} \tag{A2.16}
\end{equation*}
$$

If $f \in \mathcal{R}^{(2-n)}$, then we obviously have $J^{(2-n)} f=0$. By (A2.9) and the maximum modulus principle, we have $\operatorname{ker}\left(T^{\dagger}\right)=\mathcal{R}^{(2-n)}$. Therefore it follows from Proposition A2.3 that

$$
\int_{M}\left|\left(J^{(2-n)} f\right)(w)\right|^{2} d \mu(w)=\left\langle T^{\dagger} f, f\right\rangle_{2-n} \geq \gamma\|f\|_{n-2}^{2} \quad \text { for every } \quad f \in \mathcal{Q}^{(2-n)}
$$

Thus $J^{(2-n)}$ is an invertible operator that maps $\mathcal{Q}^{(2-n)}$ onto $\mathcal{E}$.
Obviously, $\mathcal{E}$ contains the range of the self-adjoint operator $Z$. Therefore $Z=0$ on $\mathcal{E}^{\perp}=L^{2}(M, d \mu) \ominus \mathcal{E}$. Let $\tilde{Z}$ denote the restriction of $Z$ to the invariant subspace $\mathcal{E}$. Then we have the operator decomposition

$$
\begin{equation*}
Z=\tilde{Z} \oplus 0 \tag{A2.17}
\end{equation*}
$$

corresponding to the space decomposition $L^{2}(M, d \mu)=\mathcal{E} \oplus \mathcal{E}^{\perp}$. For each $f \in \mathcal{Q}^{(2-n)}$, it is obvious that $\tilde{Z} J^{(2-n)} f=J^{(2-n)} T^{\dagger} f$. Combining this with (A2.16), we have

$$
\tilde{Z} J^{(2-n)} f=J^{(2-n)} T^{\dagger} f=J^{(2-n)}\left(J^{(2-n)}\right)^{*} J^{(2-n)} f
$$

Since $J^{(2-n)} \mathcal{Q}^{(2-n)}=\mathcal{E}$, this implies $\tilde{Z}=J^{(2-n)}\left(J^{(2-n)}\right)^{*}$. Since $J^{(2-n)}: \mathcal{Q}^{(2-n)} \rightarrow \mathcal{E}$ and $\left(J^{(2-n)}\right)^{*}: \mathcal{E} \rightarrow \mathcal{Q}^{(2-n)}$ are invertible, $\tilde{Z}$ is invertible on $\mathcal{E}$. Combining this invertibility with (A2.17), the proof of Proposition 5.6 is complete.

## Data availability

No data was used for the research described in the article.

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College of Data Science, Jiaxing University, Jiaxing 314001, China
and
Department of Mathematics, State University of New York at Buffalo, Buffalo, NY 14260, USA

E-mail: jxia@acsu.buffalo.edu

