

ORIGINAL RESEARCH ARTICLE



On the Dimension of Bergman Spaces on \mathbb{P}^1

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Received: 15 October 2021 / Revised: 6 March 2022 / Accepted: 9 March 2022 © The Author(s), under exclusive licence to Springer Science+Business Media LLC, part of Springer Nature 2022

Abstract

Inspired by a result of Szőke, we give potential-theoretic characterizations of the dimension of the Bergman space of holomorphic sections of the restriction of a holomorphic line bundle on \mathbb{P}^1 to some open set $D \subset \mathbb{P}^1$.

Keywords Bergman space · Riemann sphere · Holomorphic line bundles

Mathematics Subject Classification 32A36 · 32L05 · 30F99

1 Introduction

The Bergman space, $A^2(\Omega)$, of an open set $\Omega \subset \mathbb{C}^n$ is the vector space of $L^2(\Omega)$ -integrable holomorphic functions on Ω . Endowed with the $L^2(\Omega)$ -norm, $A^2(\Omega)$ is a reproducing kernel Hilbert space on Ω , whose reproducing kernel captures the geometry of the underlying domain. Because of that, Bergman spaces and their associated operators have been heavily studied in complex analysis. However, some fundamental properties of Bergman spaces are still unknown. For instance, open sets with infinite dimensional Bergman spaces have not been completely characterized. While this problem is interesting in its own right, it also has geometric consequences since the

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Published online: 01 April 2022

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P. Gupta was partially supported by an Infosys Young Investigator Award. L. Vivas was partially supported by National Science Foundation Grant no. 1800777.

dimension of the Bergman space is a biholomorphic invariant, and is monotonic under the inclusion of open sets. For example, these properties are used in [7] to show that a smooth strongly pseudoconvex domain in \mathbb{C}^n cannot contain a biholomorphic copy of \mathbb{C}^n , i.e., a Fatou–Bieberbach domain.

Note that holomorphic monomials are always L^2 -integrable on bounded subsets of \mathbb{C}^n . Thus, open sets that are biholomorphically equivalent to bounded sets have infinite dimensional Bergman spaces. In fact, it is known that the Bergman space of an open set in \mathbb{C} is infinite dimensional if it is non-trivial, see [18]. This dichotomy does not hold in higher dimensions, see Wiegerinck's examples [18]. We note that these examples are not pseudoconvex domains. It is an open problem whether this dichotomy holds for pseudoconvex domains, which might be viewed as true analogues of open sets in \mathbb{C} . See [3,6,7,10,12,13] for partial results on the dimension of Bergman spaces of open sets in higher dimensions.

Through work of Carleson [4] and Wiegerinck [18], the dimension of the Bergman space of an open set in the complex plane is completely characterized by the polarity of the complement, see the equivalences (a)–(c) in Theorem 1.1. The results on the dimension of the Bergman space of open sets in \mathbb{C} may be consolidated as the following theorem.

Theorem 1.1 [4], [18], [7], [8] Let $K \subsetneq \mathbb{C}$ be a closed subset. Then the following are equivalent.

- (a) K is polar.
- (b) $A^2(\mathbb{C} \setminus K) = \{0\}.$
- (c) dim $A^2(\mathbb{C} \setminus K) < \infty$.
- (d) There exists no bounded $\psi \in C^{\infty}(\mathbb{C} \setminus K)$ such that $\Delta \psi > 0$ on $\mathbb{C} \setminus K$.

That the Bergman space is either trivial or infinite dimensional, i.e., the equivalence of (b) and (c), is due to Wiegerinck, see [18]. Carleson shows in [4, Theorem 1, Sect. VI] that the Bergman space is nontrivial if and only if the logarithmic capacity of K is positive, that is, the equivalence of (a) and (c). Property (d) is an additional, potential-theoretic characterization of open sets in \mathbb{C} with finite dimensional Bergman spaces. A proof of implication $(c) \Rightarrow (d)$ may be found in [7], while implication $(d) \Rightarrow (c)$ is due to [8, Prop. 5.1].

Recently, Szőke [16] considered the dimension problem in the setting of compact Riemann surfaces. He conjectures that, given a holomorphic line bundle L on a compact Riemann surface M, the Bergman space, $A^2(D;L)$, of holomorphic sections of the restriction of L to an open set $D \subset M$ either coincides with the space of global holomorphic sections of L, or is infinite dimensional; see Sect. 2.2 for the precise definition of $A^2(D;L)$. Furthermore, he proves this conjecture in the case when M is \mathbb{P}^1 , the Riemann sphere.

Theorem 1.2 [16] Let L be a holomorphic line bundle on \mathbb{P}^1 , and $D \subset \mathbb{P}^1$ be an open set. Then $A^2(D; L)$ is either equal to $\Gamma(\mathbb{P}^1; L)$, i.e., the space of global holomorphic sections on \mathbb{P}^1 , or infinite dimensional.

Szőke's proof is a modification of Wiegerinck's proof in [18], and relies on the relatively simple algebraic structure of holomorphic line bundles on \mathbb{P}^1 . While this proof



is hard to generalize to other compact Riemann surfaces, Szőke gives a partial result for general compact Riemann surfaces. However, he does not provide any potential-theoretic characterizations, akin to Theorem 1.1, for the general case; see Proposition 2.1 in [16]. An objective of this paper is to fill this gap for \mathbb{P}^1 . In fact, our main result is a complete analogue of Theorem 1.1 for Bergman spaces of holomorphic sections on open subsets of \mathbb{P}^1 .

Theorem 1.3 Let $\Pi: L \to \mathbb{P}^1$ be a holomorphic line bundle. Suppose that $K \subsetneq \mathbb{P}^1$ is a compact subset. Then the following are equivalent.

- (a) K is polar.
- (b) $A^2(\mathbb{P}^1 \setminus K; L) = \Gamma(\mathbb{P}^1; L)$.
- (c) dim $A^2(\mathbb{P}^1 \setminus K; L) < \infty$.
- (d) There exists no bounded function $\psi \in C^{\infty}(\mathbb{P}^1 \setminus K)$ such that $i \partial \bar{\partial} \psi \geq \omega$ on $\mathbb{P}^1 \setminus K$ for some volume form ω on \mathbb{P}^1 .

The implications $(a) \Rightarrow (b) \Rightarrow (c) \Rightarrow (a)$ are proved using some basic results in potential theory. In particular, our proof of $(c) \Rightarrow (a)$ in Sect. 3.2 is inspired by Wiegerinck's proof in [18], but uses Carleson's construction of a nontrivial L^2 -integrable holomorphic function as seen in [4, Theorem 9.5, Chap. 21]. We include this proof of the implication $(c) \Rightarrow (a)$, because it does not rely on partial differential equations techniques, unlike our proof of $(c) \Rightarrow (d) \Rightarrow (a)$.

As in the proof of Theorem 1.1, the proof of the implication $(c) \Rightarrow (d)$ in Theorem 1.3 is done by using Hörmander's weighted L^2 -method for solving the Cauchy–Riemann equations. We note that the equivalence of (a) and (d) can be paraphrased as: K is nonpolar if and only if, for every Hermitian metric h on L and volume form ω on \mathbb{P}^1 , there is a bounded function $\psi \in \mathcal{C}^{\infty}(\mathbb{P}^1 \setminus K)$ satisfying

$$i\Theta_{a^{-\psi}h} > \omega \quad \text{on } \mathbb{P}^1 \setminus K.$$

That is, whenever K is nonpolar, the given Hermitian metric h on L may be twisted such that the newly obtained Hermitian metric on $L|_D$ has positive curvature, see [11, Theorem 3.11.2] for a related result on general open Riemann surfaces. This positivity lets one use Hörmander's method successfully; the boundedness of ψ ensures that the resultant estimates are for the Bergman spaces associated to the given Hermitian metric h on L.

The proof of $(d) \Rightarrow (a)$ in Theorem 1.3 is also motivated by the proof of the same implication in Theorem 1.1. In the latter, one uses the potential function associated to the equilibrium measure of a compact, nonpolar set to construct the strictly subharmonic, bounded weight function ψ . This construction may be done in a local chart to prove the implication $(d) \Rightarrow (a)$ in Theorem 1.3. In fact, the proofs of the equivalences in Theorem 1.3 are all done in a local chart. However, these proofs, except for the direct proof of $(c) \Rightarrow (a)$ in Sect. 3.2, do not depend on the particular structure of L and h so that an extension of Theorem 1.3 to compact Riemann surfaces seems feasible. In fact, we will consider this question for compact Riemann surfaces in a forthcoming paper.

As a byproduct of the proof of the implication $(c) \Rightarrow (d)$ of Theorem 1.3, one obtains a complete description of the dimension of the weighted Bergman space



 $A_{e^{-\psi}}^2(\Omega)$, where $\Omega\subset\mathbb{C}$ is an open set, and ψ is a subharmonic function on \mathbb{C} . This is a combination of Corollary 3.3, which is a version of Theorem 1.1 for $A_{e^{-\psi}}^2(\Omega)$, and the complete characterization of the dimension of $A_{e^{-\psi}}(\mathbb{C})$ due to Borichev, Le, and Youssfi; see [2, Theorem 2.6].

In [18], Wiegerinck gives examples of domains in \mathbb{C}^2 that have nontrivial, but finite dimensional Bergman spaces. These examples immediately show that the dichotomy conjectured for compact Riemann surfaces does not hold for the Bergman spaces of holomorphic sections of a specific holomorphic line bundle, $L = \mathcal{O}(-3)$, on the complex projective space \mathbb{P}^2 . Motivated by Wiegerinck's construction, we produce non-pseudoconvex examples in \mathbb{P}^2 to show that this dichotomy is absent irrespective of the choice of holomorphic line bundle on \mathbb{P}^2 .

The paper is structured as follows. In Sect. 2, we detail the required background material and notation on Bergman spaces, potential theory, and holomorphic line bundles on \mathbb{P}^1 . The first four subsections of Sect. 3 contain the proofs of the equivalences of Theorem 1.3. This is followed by Sect. 3.5, where the results on the dimension of weighted Bergman spaces for open sets in \mathbb{C} are given. In Sect. 4, we show by example that for each holomorphic line bundle on \mathbb{P}^2 , there is a domain in \mathbb{P}^2 such that the Bergman space of the corresponding holomorphic sections does not equal the space of global holomorphic sections but is finite dimensional.

2 Background and Preliminaries

2.1 Bergman Spaces in \mathbb{C}^n and Potential Theory in \mathbb{C}

Let $\Omega \subset \mathbb{C}^n$ be an open subset, and λ be the standard volume form on \mathbb{C} . For a positive function ϕ on Ω , the *weighted Bergman space* of Ω with weight ϕ is the Banach space

$$A_{\phi}^2(\Omega) := \left\{ f \in \mathcal{O}(\Omega) : ||f|| := \left(\int_{\Omega} |f(z)|^2 \phi(z) \lambda(z) \right)^{1/2} < \infty \right\}.$$

For $\phi \equiv 1$, the space $A_{\phi}^2(\Omega)$ is denoted by $A^2(\Omega)$, and referred to as the Bergman space of Ω .

A set $K \subset \mathbb{C}$ is said to be *polar* if there is a nonconstant subharmonic function s on \mathbb{C} such that $K \subset \{z \in \mathbb{C} : s(z) = -\infty\}$. Polar sets admit an alternate characterization via logarithmic potential theory. For a finite Borel measure ν with compact support in \mathbb{C} , its *potential* is the function $p_{\nu} : \mathbb{C} \to [-\infty, \infty)$ given by

$$p_{\nu}(z) = \int_{\mathbb{C}} \ln|z - w| \, d\nu(w), \qquad z \in \mathbb{C}.$$

The *energy* of ν is the quantity

$$I(v) = \int_{\mathbb{C}} p_{v}(z) \; dv(z).$$

The (logarithmic) capacity of $K \subset \mathbb{C}$ is defined as

 $\operatorname{cap}(K) := \sup\{e^{I(\nu)} : \nu \text{ is a Borel probability measure with compact support in } K\}.$



For a compact nonpolar set $K \subset \mathbb{C}$, it is known that there is a unique Borel probability measure ν_K supported on K such that

$$I(\nu_K) = \sup\{I(\nu) : \nu \text{ is a Borel probability measure on } K\}.$$

This measure is called the *equilibrium measure* of *K*. The following result relates the notion of polarity, local polarity, and capacity.

Theorem 2.1 *Let* $K \subset \mathbb{C}$ *be a Borel subset. Then the following are equivalent.*

- (i) K is polar.
- (ii) There is an open set G containing K, and a subharmonic function s on G, nonconstant on any component of G, such that $K \subset \{x \in G : s(x) = -\infty\}$.
- (iii) For any $\zeta \in \mathbb{C}$, there is a connected open neighborhood V_{ζ} of ζ and a nonconstant subharmonic function s_{ζ} on V_{ζ} such that $K \cap V_{\zeta} \subset \{x \in V_{\zeta} : s_{\zeta}(x) = -\infty\}$.
- (iv) cap(K) = 0.

For the equivalences of (i) and (ii), (i) and (iii), and (i) and (iv), see for instance Proposition 5.5, Lemma 5.6 and Theorem 7.5, respectively, in [5, Chap. 21].

2.2 Bergman Spaces of Holomorphic Sections

Let M be a complex manifold. Given a holomorphic line bundle $\Pi: L \to M$ and an open set $D \subset M$, the space of holomorphic sections of $L|_D$ is denoted by $\Gamma(D; L)$.

Let ω be a volume form on M. Given a holomorphic line bundle $\Pi: L \to M$, and a smooth Hermitian metric h on L, we define the Bergman space of sections of $(L|_D, h)$ as

$$A^2_{h,\omega}(D;L) = \left\{ s \in \Gamma(D;L) : ||s|| := \left(\int_D h(s,s) \, \omega \right)^{1/2} < \infty \right\}.$$

 $A^2_{h,\omega}(D;L)$ is a reproducing kernel Hilbert space. If $M=\mathbb{C}^n$, then $\Gamma(D;L)=\mathcal{O}(D)$ and $h(s,s)\omega=|s|^2\phi\lambda$, for some positive function ϕ on \mathbb{C}^n . In this case, $A^2_{h,\omega}(D;L)$ is the weighted Bergman space $A^2_{\phi}(D)$.

Lemma 2.2 Let M, D, L, h, ω be as above. Suppose that M is compact. Then, as a vector space, $A_{h,\omega}^2(D;L)$ is independent of the choices of h and ω . In particular, $A_{h,\omega}^2(M;L) = \Gamma(M;L)$.

Proof Given any two smooth volume forms, ω and ω' , on M, and any two smooth Hermitian metrics, h and h', on L, there exist smooth positive functions, f and g, on M, such that $\omega = f\omega'$ and h = gh'. Thus, the Bergman spaces of the sections of $L|_D$ are isomorphic for different choices of the volume form on M and Hermitian metric on L.

The final claim holds since h(s, s) is a continuous positive function on M for any choice of h and s.



Remark 2.3 In this paper, M is either \mathbb{P}^1 or \mathbb{P}^2 . Thus, in view of Lemma 2.2, $A_{h,\omega}^2(D;L)$ is simply denoted by $A^2(D;L)$.

2.3 Line Bundles on \mathbb{P}^1

We recall certain standard facts about \mathbb{P}^1 . Given any $q \in \mathbb{P}^1$, one may choose coordinates $\zeta = [\zeta_0 : \zeta_1]$ on \mathbb{P}^1 so that q = [0 : 1]. In this case, we often identify $U^z = \mathbb{P}^1 \setminus \{q\}$ with the complex plane \mathbb{C} via

$$\varphi_z: \zeta = [\zeta_0: \zeta_1] \mapsto z = \frac{\zeta_1}{\zeta_0},$$

and refer to z as the local coordinate on the affine chart $\mathbb{P}^1 \setminus \{q\}$. On occasion, the chart $(U^{1/z}, \varphi_{1/z})$ given by $\varphi_{1/z} : \zeta \mapsto \zeta_0/\zeta_1$ on $U^{1/z} = \mathbb{P}^1 \setminus \{[1:0]\}$ will also be used. In this case, we will use the coordinate w on $\varphi_{1/z}(U^{1/z}) = \mathbb{C}$.

A smooth volume form ω on \mathbb{P}^1 is of the form f ω_{FS} , where ω_{FS} is the Fubini–Study volume form, and f is a smooth positive function on \mathbb{P}^1 . Thus, in the local coordinate z,

$$\omega(z) = f(z) \frac{idz \wedge d\overline{z}}{2(1+|z|^2)^2}, \quad z \in \mathbb{C},$$

where $f:\mathbb{C}\to (0,\infty)$ is smooth, and $f\circ\varphi_z$ admits a smooth positive extension to \mathbb{P}^1 . Given a smooth function ψ on an open set in \mathbb{P}^1 , $\partial\overline{\partial}\psi$ is the unique smooth (1, 1)-form on \mathbb{P}^1 satisfying $(\varphi_z)_*\left(\partial\overline{\partial}\psi\right)(z)=\frac{\partial^2(\psi\circ\varphi_z)}{\partial z\partial\overline{z}}(z)\,dz\wedge d\overline{z}$ for all $z\in U^z$. Next, any holomorphic line bundle $\Pi:L\to\mathbb{P}^1$ is of the form $\mathcal{O}(k)$, for some $k\in\mathbb{Z}$,

Next, any holomorphic line bundle $\Pi: L \to \mathbb{P}^1$ is of the form $\mathcal{O}(k)$, for some $k \in \mathbb{Z}$, where $\mathcal{O}(k)$ is the line bundle associated to the divisor $k\{p\}$ for any fixed $p \in \mathbb{P}^1$. In the local coordinate z, a global section of $\mathcal{O}(k)$ is given by a pair, $s = (s_1, s_2)$, of holomorphic functions on \mathbb{C} such that $s_1(z) = z^k s_2(1/z)$ if $z \in \mathbb{C}^*$. Similarly, a Hermitian metric on $L = \mathcal{O}(k)$ is given by a pair, $h = (h_1, h_2)$, of smooth positive functions on \mathbb{C} , such that $h_1(z) = |z|^{-2k}h_2(1/z)$ if $z \in \mathbb{C}^*$. Owing to the relationship between h_1 and h_2 , $\partial \overline{\partial} h_1(z) = \partial \overline{\partial} h_2(1/z)$, $z \in \mathbb{C}^*$, where $\partial \overline{\partial} f(z) = \frac{\partial^2 f}{\partial z \partial \overline{z}}(z) \, dz \wedge d\overline{z}$ for any \mathcal{C}^2 -smooth function f on \mathbb{C} . Thus, there is a smooth (1, 1)-form Θ_h on \mathbb{P}^1 such that $(\varphi_z)_*\Theta_h = \partial \overline{\partial} h_1$ and $(\varphi_{1/z})_*\Theta_h = \partial \overline{\partial} h_2$. The form Θ_h is referred to as the curvature of the Hermitian metric h. Note that it suffices to specify s_1 , h_1 and $\partial \overline{\partial} h_1$, which is the convention we will employ. Lastly, recall the following description of the space of global sections of $\mathcal{O}(k)$, $k \in \mathbb{Z}$.

Lemma 2.4 *Let* $k \in \mathbb{Z}$. *In the local coordinate* z,

$$A^{2}(\mathbb{P}^{1}; \mathcal{O}(k)) = \Gamma(\mathbb{P}^{1}; \mathcal{O}(k)) = \begin{cases} \{polynomials \ of \ degree \ at \ most \ k\}, & if \ k \geq 0, \\ \{f \equiv 0\}, & if \ k < 0. \end{cases}$$

In particular, dim $A^2(\mathbb{P}^1; \mathcal{O}(k)) = \dim \Gamma(\mathbb{P}^1; \mathcal{O}(k)) = \max\{0, k+1\}.$

While Lemma 2.4 is a standard result, see [17, p. 81], we direct the reader to the remarks before Corollary 3.3 for a proof in the case $k \ge -2$.



2.4 Potential Theory on \mathbb{P}^1

Given an open set $U \subset \mathbb{P}^1$, and an upper semicontinuous function $s: U \to [-\infty, \infty)$ with $s \not\equiv -\infty$ on any connected component of U, s is said to be *subharmonic on* U if, for every coordinate chart (V, φ) that intersects $U, s \circ \varphi^{-1}$ is subharmonic on $\varphi(U \cap V) \subset \mathbb{C}$.

Definition 2.5 Given a set $K \subset \mathbb{P}^1$, K is said to be *polar* if for each $\zeta \in \mathbb{P}^1$, there is an open neighborhood $V_{\zeta} \subset \mathbb{P}^1$ of ζ and a subharmonic function s_{ζ} on V_{ζ} such that $K \cap V_{\zeta} \subset \{x \in V_{\zeta} : s_{\zeta}(x) = -\infty\}$.

Lemma 2.6 Let $K \subset \mathbb{P}^1$ be compact. For any fixed $q = [0:1] \in \mathbb{P}^1$, K is polar in \mathbb{P}^1 if and only if $K_* := \varphi_7(K)$ is polar in \mathbb{C} .

Proof First, suppose that K is polar. Let $\zeta_* \in \mathbb{C}$ and $\zeta = \varphi_z^{-1}(\zeta_*)$. Then it is clear that

$$K_* \cap U_{\zeta_*} \subset \{x \in U_{\zeta_*} : s_{\zeta_*}(x) = -\infty\},\$$

where $U_{\zeta_*} = \varphi_z(V_{\zeta} \setminus \{q\})$, and $s_{\zeta_*} = s_{\zeta} \circ \varphi_z^{-1}$. Thus, K_* is polar in \mathbb{C} .

Next, assume that K_* is polar. By a similar reasoning as above, we obtain that for every $\zeta \in \mathbb{P}^1 \setminus \{q\}$, there exists a $V_\zeta \subset \mathbb{P}^1$ and a subharmonic function s_ζ on V_ζ that satisfy the condition in Definition 2.5. It remains to produce such a pair (V_ζ, s_ζ) for $\zeta = q$. Since subharmonicity is preserved by biholomorphic maps, it follows that the set $K_{**} := \varphi_{1/z} \circ \varphi_z^{-1}(K_*)$ is polar in \mathbb{C}^* . In fact, by Theorem 2.1, K_{**} is polar in \mathbb{C} . On the other hand, $\{0\} = \varphi_{1/z}(\{q\})$ is also polar in \mathbb{C} . Now, using the fact that the union of two polar sets in \mathbb{C} is polar in \mathbb{C} , see [5, Lemma 5.6, Chap. 21], $\varphi_{1/z}(K) = K_{**} \cup \{0\}$ is polar in \mathbb{C} , i.e., there is a subharmonic function s on \mathbb{C} such that $\varphi_{1/z}(K) \subset \{x \in \mathbb{C} : s(x) = -\infty\}$. To complete the proof, we simply observe that $V_q = \varphi_{1/z}^{-1}(\mathbb{C})$ and $s_q = s \circ \varphi_{1/z}$ satisfy the condition required at $q \in K$ in Definition 2.5.

3 Proof(s) of the Main Result

By Lemma 2.4, $(b) \Rightarrow (c)$. In Sects. 3.1–3.4, we prove the implications $(a) \Rightarrow (b)$, $(c) \Rightarrow (a)$, $(c) \Rightarrow (d)$, and $(d) \Rightarrow (a)$, respectively. This section concludes with the characterization of weighted Bergman spaces, $A_{e^{-\psi}}^2(\Omega)$, for subharmonic ψ and $\Omega \subset \mathbb{C}$ open.

3.1 Proof of (a) Implies (b)

In order to prove that $A^2(\mathbb{P}^1 \setminus K; L) = \Gamma(\mathbb{P}^1; L)$ whenever $K \subset \mathbb{P}^1$ is a compact polar set, we first prove a version of this result for sets in \mathbb{C} in Lemma 3.1 below. A readily available result in the literature is that $A^2(\Omega) = A^2(\Omega \setminus X)$ for any open set $\Omega \subset \mathbb{C}$ and polar, compact set $X \subset \Omega$, see, for instance, Theorem 9.5 in [5, Chap. 21]



and the paragraph preceding it. We present a generalization of this result to the case when X is a polar and relatively closed subset of Ω .

Lemma 3.1 [15] Let $\Omega \subset \mathbb{C}$ be an open subset, and $X \subset \mathbb{C}$ be a closed polar subset. Then, $A^2(\Omega \setminus X) = A^2(\Omega)$, i.e., for any $f \in A^2(\Omega \setminus X)$, there is an $F \in A^2(\Omega)$ such that $F|_{\Omega \setminus X} = f$.

Proof The proof uses the fact that for any $a \in \Omega \cap X$, there exists an $r_a > 0$ such that $\overline{\mathbb{D}(a, r_a)} \subset \Omega$, and $b\mathbb{D}(a, r_a) \cap X = \emptyset$, see Theorem 7.3.9 in [1].

Now, let $D_a:=\mathbb{D}(a,r_a)$ and $X_a:=\mathbb{D}(a,r_a)\cap X$, where r_a is as above. Then, X_a is a compact subset of D_a . Using the result for compact, polar sets mentioned before the statement of this lemma, for any $f\in A^2(\Omega)$, there is an $F_a\in A^2(D_a)$ such that $F_a|_{D_a\setminus X_a}=f|_{D_a\setminus X_a}$. For any $a,b\in X$ such that $D_a\cap D_b\neq\emptyset$, we have that $(D_a\cap D_b)\setminus X$ is a nonempty open set because X is a closed and polar set. Thus, it follows from the identity theorem that $F_a=F_b$ on $D_a\cap D_b$. Hence, the following function is well-defined and holomorphic on Ω :

$$F(z) = \begin{cases} F_a(z), & z \in D_a, \\ f(z), & z \in \Omega \setminus X. \end{cases}$$

Moreover, since f and F differ only on the zero measure set X, it follows that $||F||_{A^2(\Omega)} = ||f||_{A^2(\Omega \setminus X)}$.

We are now set to prove that (a) implies (b) in Theorem 1.3.

Proof (Proof of $(a) \Rightarrow (b)$ in Theorem 1.3) Suppose that $K \subsetneq \mathbb{P}^1$ is polar. If K is empty, then (b) follows from (a), see Lemma 2.4. Thus, we assume that K is a nonempty polar set, and fix a $q \in K$. With q = [0:1], we fix the Hermitian metric on $L = \mathcal{O}(k)$ as $h_1(z) = 1/(1+|z|^2)^k$ in the local coordinate z. To account for the volume form in the local coordinate z, we set

$$\phi_k(z) := \frac{1}{(1+|z|^2)^{k+2}}, \qquad z \in \mathbb{C}. \tag{3.1}$$

Setting $D_* := \mathbb{C} \setminus K_*$, it then follows that $A^2(D; \mathcal{O}(k))$ is isomorphic to $A^2_{\phi_k}(D_*)$.

Next, we show that $A_{\phi_k}^2(D_*) = A_{\phi_k}^2(\mathbb{C})$. For this, let $f \in A_{\phi_k}^2(D_*)$. For any R > 0, set $D_*(R) = D_* \cap \mathbb{D}(0; R)$ and $f_R = f \cdot \chi_{D_*(R)}$. Then, since

$$\phi_k(z) \geq \begin{cases} 1, & \text{if } k < -2, \\ (1+R^2)^{-k-2}, & \text{if } k \geq -2, \end{cases} \quad z \in \mathbb{D}(0; R),$$

it follows that $f_R \in A^2(D_*(R))$. However, K_* is a closed polar subset of $\mathbb C$. Thus, by Lemma 3.1, there exists an $F_R \in A^2(\mathbb D(0;R))$ such that $F_R|_{D_*(R)} = f_R$. We abuse

¹ After the first draft of this manuscript appeared on the arXiv, Pflug brought to our attention that Lemma 3.1 has appeared in [15], where Siciak attributes it to Sakai and Skwarczyński. Since neither [15] nor the source cited therein are easily accessible, we give a proof of Lemma 3.1 in this work. This proof is simpler than our original one, and follows an argument suggested by Pflug. This argument resembles that of Siciak's.



notation, and assume that F_R is defined on \mathbb{C} by setting it to be 0 on $\mathbb{C} \setminus \mathbb{D}(0; R)$. Since

$$F_S|_{D_*(R)} = F_R|_{D_*(R)} = f|_{D_*(R)}$$
 whenever $S \ge R > 0$,

it follows from the identity theorem that $F_S|_{\mathbb{D}(0;R)} = F_R|_{\mathbb{D}(0;R)}$ for all $S \geq R > 0$. Thus, the sequence $\{F_N\}_{n \in \mathbb{N}}$ admits a pointwise limit, say F, on \mathbb{C} . Moreover, since $F|_{\mathbb{D}(0;N)} = F_N$ for all $N \in \mathbb{N}$, $F \in \mathcal{O}(\mathbb{C})$ and $F|_{D_*} = f$. In particular, F is measurable so that

$$\|F\|_{L^2_{\phi_k}(\mathbb{C})} = \|F\|_{L^2_{\phi_k}(D_*)} = \|f\|_{L^2_{\phi_k}(D_*)} < \infty$$

follows. Therefore, $A_{\phi_k}^2(D_*)$ is isomorphic to $A_{\phi_k}^2(\mathbb{C})$. Finally, note that $A_{\phi_k}^2(\mathbb{C})$ is isomorphic to $A^2(\mathbb{P}^1; \mathcal{O}(k))$, which is $\Gamma(\mathbb{P}^1; \mathcal{O}(k))$, by Lemma 2.4. Thus, (a) implies (b) in Theorem 1.3.

3.2 Proof of (c) Implies (a) via the Cauchy Transform Approach

Proof We prove the implication by contraposition, i.e., we prove that if $K \subsetneq \mathbb{P}^1$ is nonpolar, then $A^2(\mathbb{P}^1 \setminus K; \mathcal{O}(k))$ is infinite dimensional. As in the proof of $(a) \Rightarrow (b)$ in Sect. 3.1, we assume that $q = [0:1] \in K$ and fix the Hermitian metric $h_1(z) = 1/(1+|z|^2)^k$ on $\mathcal{O}(k)$. Then, it suffices to show that $A^2_{\phi_k}(D_*)$ is infinite dimensional, where ϕ_k is as in (3.1), and $D_* = \mathbb{C} \setminus K_*$.

Case 1. Suppose that $k \geq -2$. Then, since $\phi_k(z) \leq 1$ for all $z \in \mathbb{C}$, $A^2(D_*) \subset A^2_{\phi_k}(D_*)$. Since, by Lemma 2.6, $\mathbb{C} \setminus D_* = K_*$ is nonpolar, Theorem 1.1 yields that $A^2(D_*)$, and therefore, $A^2_{\phi_k}(D_*)$ is infinite dimensional.

Case 2. Suppose that k < -2. In this case, $A_{\phi_k}^2(D_*) \subset A^2(D_*)$, so Theorem 1.1 does not directly yield our claim. However, we use the techniques of Carleson and Wiegerinck to produce infinitely many independent functions in $A_{\phi_k}^2(D_*)$.

First, we recall Carleson's construction of a nontrivial function in $A^2(D_*)$; see [5, Theorem 9.5, Chap. 21] for a detailed exposition. Set $K_* := \varphi_z(K) = \mathbb{C} \setminus D_*$. Then, K_* is a Borel nonpolar set in \mathbb{C} . Thus, it contains a compact set, say E, with positive logarithmic capacity, see for instance [5, Theorem 7.5, Chap. 21]. Let E_1 , E_2 be disjoint compact subsets of E, each of which has positive logarithmic capacity. For $j \in \{1, 2\}$, let μ_j be the equilibrium measure of E_j , and set $\mu = \mu_1 - \mu_2$. Set f to be the Cauchy transform of μ , i.e.,

$$f(z) = \int_{F} \frac{d\mu(\xi)}{\xi - z}.$$
 (3.2)

Then, f is analytic on $\mathbb{C}_{\infty} \setminus E$, with $f(\infty) = f'(\infty) = 0$. Carleson further shows that $f \in A^2(\mathbb{C} \setminus E)$ which is contained in $A^2(D_*)$.

We next use Wiegerinck's technique to produce a sequence of linearly independent functions $\{g_j\}_{j\in\mathbb{N}}\subset A^2(\mathbb{C}\setminus E)$ such that g_j vanishes at ∞ , and its order of vanishing



at ∞ is at least j. Assuming the existence of this sequence for the moment, we claim that $g_j \in A^2_{\phi_k}(D_*)$ for all $j \geq -k$. To see this, fix a $j \geq -k$, and note that there exist $\{c_{j,\ell}\}_{\ell \geq j} \subset \mathbb{C}$ and R > 0 so that

$$g^{j}(z) = \sum_{\ell=j}^{\infty} c_{\ell,j} z^{-\ell} \quad \text{for } |z| > R.$$

Thus, since $j \ge -k$ and k < -2, it follows that

$$\begin{split} \int_{D_*} |g_j(z)|^2 \frac{idz \wedge d\overline{z}}{2(1+|z|^2)^{k+2}} \\ & \leq (1+R^2)^{-2-k} \int_{D(0;R)\backslash E} |g_j(z)|^2 \, \lambda(z) \\ & + \int_{|z|>R} |g_j(z)|^2 \frac{idz \wedge d\overline{z}}{2(1+|z|^2)^{k+2}} \\ & \leq ||g_j||_{A^2(\mathbb{C}\backslash E)}^2 + 2\pi \int_{r>R} \sum_{\ell=j} |c_{\ell,j}|^2 r^{-2\ell} r (1+r^2)^{-2-k} dr < \infty. \end{split}$$

This gives the infinite dimensionality of $A_{\phi_k}^2(D_*)$. It remains to construct the sequence $\{g_j\}_{j\in\mathbb{N}}$.

Subcase 1. Suppose that f in (3.2) is rational. Since f is L^2 -integrable on $\mathbb{C} \setminus E$, but not on \mathbb{C} , E must have positive Lebesgue measure. In this case, the function

$$g(z) = \int_{F} \frac{\lambda(\xi)}{\xi - z}$$

is bounded and analytic on $\mathbb{C} \setminus E$, with $g(\infty) = 0$ and $g'(\infty) = -\lambda(E)$, see [9, p. 2]. Thus, $g^j \in A^2(\mathbb{C} \setminus E)$ for all $j \in \mathbb{N}$, and each g^j has order of vanishing j at ∞ . In this case, we set $g_j = g^j$.

Subcase 2. Suppose that f in (3.2) is not rational. Expanding f in a Laurent series around ∞ , one gets

$$f(z) = \sum_{\ell=p}^{\infty} c_{\ell} z^{-\ell}$$
, for some $p \ge 2$, $c_p \ne 0$.

Now, we produce a nontrivial function $g \in A^2(\mathbb{C} \setminus E)$ whose Laurent expansion at ∞ does not contain any terms in $z^{-1},...,z^{-p}$. Let $z_1,...,z_{p+1}$ be distinct points in $\mathbb{C} \setminus E$, and

$$g(z) = \sum_{\ell=1}^{p+1} b_{\ell} \frac{(f(z) - f(z_{\ell}))}{z - z_{\ell}}.$$

Then, expanding g as a Laurent series around ∞ , one obtains $g(z) = \sum_{\ell=1}^{\infty} a_{\ell} z^{\ell}$, where

$$a_m = \sum_{\ell=1}^{p+1} -b_\ell f(z_\ell) z_\ell^{m-1}, \qquad m \in \{1, \dots, p\}.$$
 (3.3)



Here, the constants b_1, \ldots, b_{p+1} are chosen to solve the homogeneous system of p linear equations obtained by setting $a_m = 0, m \in \{1, \ldots, p\}$. Moreover, g cannot be trivial, else f will be rational. Thus, g has the desired properties.

Setting $g_1 = \cdots = g_p = f$ and $g_{p+1} = g$, we construct g_j inductively for $j \ge p+1$ by repeating the above procedure for g_{j-1} in place of f. This completes the construction of the sequence in all cases, and hence, the proof of Theorem 1.3. \square

3.3 Proof of (c) Implies (d)

The proof of this implication follows from a similar result for sets in \mathbb{C} .

Lemma 3.2 Let $\Omega \subset \mathbb{C}$ be an open set and $A^2_{e^{-\psi_1}}(\Omega)$ be the weighted Bergman space for some $\psi_1 \in \mathcal{C}^{\infty}(\Omega, \mathbb{R})$. Suppose that there exists a subharmonic function $\psi_2 \in \mathcal{C}^{\infty}(\Omega, \mathbb{R})$ such that it is bounded above, $\psi_1 + \psi_2$ is subharmonic on Ω , and

$$\Delta \psi_2 > 0$$
 on \overline{U}

for an open set $U \subseteq \Omega$. Then $A^2_{e^{-\psi_1}}(\Omega)$ is infinite dimensionial.

Proof This follows from work in [7], see Theorem 1 and Lemma 7 therein. The latter requires a slight modification. This modification consists of replacing $K\varphi(z)$ by $\psi_1(z) + K\psi_2(z)$, and using the boundedness from above of ψ_2 to show that the constructed function u belongs to $L^2_{a-\psi_1}(\Omega)$.

Proof (Proof of $(c) \Rightarrow (d)$ in Theorem 1.3) The proof is done by contraposition. Thus, we assume that there exists a bounded function $\psi \in \mathcal{C}^{\infty}(\mathbb{P}^1 \setminus K)$ such that $i \partial \bar{\partial} \psi \geq \omega$ holds on $\mathbb{P}^1 \setminus K$ for some volume form ω on \mathbb{P}^1 . Note that, if K was empty, then it would follow that there is a bounded, nonconstant subharmonic function on \mathbb{C} . Hence, we may assume that K is nonempty. As before, we choose a $q \in K$ and coordinates such that q = [0:1]. Set $K_* = \varphi_Z(K)$.

Next, we note that, as in Sect. 3.2, it suffices to show that $A_{\phi_k}^2(D_*)$, for $D_* := \mathbb{C} \setminus K_*$, is infinite dimensional for each $k \in \mathbb{Z}$.

For that, we use Lemma 3.2. In particular, we set $\widetilde{\psi}_2 = \psi \circ \varphi_z^{-1}$. Then, clearly, $\widetilde{\psi}_2 \in \mathcal{C}^{\infty}(D_*)$ is a bounded, strictly subharmonic function. Moreover, it follows that there is a constant $c_1 > 0$ such that

$$\Delta \widetilde{\psi}_2(z) \ge c_1 (1 + |z|^2)^{-2}$$
 for $z \in D_*$.

Next, set $\psi_1 := -\ln(\phi_k)$. Then $\psi_1 \in \mathcal{C}^{\infty}(D_*)$ and

$$\Delta \psi_1(z) = 4(k+2)(1+|z|^2)^{-2} \tag{3.4}$$

on D_* . Hence, there exists a constant c > 0 such that

$$\Delta(\psi_1 + c\widetilde{\psi}_2) > 0$$
 on D_* .



Setting $\psi_2=c\widetilde{\psi}_2$, it follows that the hypotheses of Lemma 3.2 are satisfied for any open set $U \in D_*$. Thus, $A^2_{\phi_k}(D_*)$ is infinite dimensional.

3.4 Proof of (d) Implies (a)

Proof This proof is also done by contraposition. That is, we assume that $K \subset \mathbb{P}^1$ is nonpolar. Since K is nonempty, we may construct ψ in a chart first. As before, we let $q \in K$ and choose coordinates such that q = [0:1]. We set $K_* = \varphi_z(K)$ and $D_* = \mathbb{C} \setminus K_*$. We shall use a function which was initially constructed in the proof of Proposition 5.1 in [8]. For that, let $G \subset K_*$ be a nonpolar, compact set and $v = v_G$ be the equilibrium measure of G. The associated potential function $p = p_G : \mathbb{C} \longrightarrow [-\infty, \infty)$ is defined as

$$p(z) = \int_{\mathbb{C}} \ln|z - w| \ d\nu(w) \text{ for } z \in \mathbb{C}.$$

Since p is harmonic on G^c , see for instance [14, Theorem 3.1.2], it follows that $\Delta e^{-p} = e^{-p} |\nabla p|^2$ on G^c . Thus, e^{-p} is subharmonic on G^c , and strictly subharmonic at all points in G^c at which the gradient of p is nonvanishing. It follows from the proof of Proposition 5.1 in [8] that $|\nabla p|$ is strictly positive outside a sufficiently large disc containing G. In particular, there exist constants τ_1 , R > 0 such that $G \subset \mathbb{D}(0, R)$ and

$$|\nabla p(z)| > \frac{\tau_1}{|z|} \text{ for } z \in \mathbb{D}(0, 2R)^c.$$

Furthermore, since ν is an equilibrium measure, one obtains for $z \in \mathbb{D}(0, 2R)^c$ that

$$p(z) = \int_{\mathbb{C}} \ln|z - w| \ d\nu(w) \le \ln(3|z|/2).$$

Therefore,

$$\triangle e^{-p(z)} \ge \frac{2}{3} \frac{\tau_1^2}{|z|^3} \text{ for } z \in \mathbb{D}(0, 2R)^c.$$

Strict subharmonicity on all of G^c may now be achieved by adding to e^{-p} a particular compactly supported function which is strictly subharmonic on $\mathbb{D}(0, 2R)$. For instance, let $\chi \in \mathcal{C}_0^\infty(\mathbb{D}(0, R'))$ for some R' > 2R such that $\chi(z) = |z|^2$ on $\mathbb{D}(0, 2R)$. Then, for $\epsilon > 0$ sufficiently small, there exist $\tau_2, \tau_3 > 0$ such that the function $\psi_*(z) := e^{-p(z)} + \epsilon \chi(z)$ satisfies

$$\Delta \psi_*(z) \ge \begin{cases} \frac{\tau_2}{|z|^3}, & \text{for } z \in \mathbb{D}(0, 2R)^c, \\ \tau_3, & \text{for } z \in \mathbb{D}(0, 2R). \end{cases}$$
(3.5)



It follows that $\psi := \psi_* \circ \varphi_7$ satisfies $i \partial \bar{\partial} \psi \geq \omega$ on $\mathbb{P}^1 \setminus K$ for some volume form ω on \mathbb{P}^1 . Moreover, Frostman's theorem implies that e^{-p} is bounded on G^c , and hence, ψ_* is bounded on G^c . That is, ψ is bounded on $\mathbb{P}^1 \setminus K$, and hence satisfies (d) in Theorem 1.3.

3.5 A Corollary to Lemma 3.2 on Weighted Bergman Spaces

A variation of Lemma 3.2 yields an analogue of Theorem 1.1 for the Bergman space $A_{e^{-\psi}}^2(\Omega)$ for ψ subharmonic on Ω , see Corollary 3.3 below.

The dimension of $A^2_{e^{-\psi}}(\mathbb{C})$ for ψ subharmonic on \mathbb{C} is completely described by Borichev, Le, and Youssfi in [2, Theorem 2.6] in terms of the continuous part $(\mu_{\psi})^c$ of the Riesz measure μ_{ψ} of ψ . They show that dim $A_{a-\psi}^2(\mathbb{C}) < \infty$ is finite if and only if $(\mu_{\psi})^c(\mathbb{C}) < \infty$, and

- (a) if $(\mu_{\psi})^c(\mathbb{C}) = 0$, then dim $A_{\rho-\psi}^2(\mathbb{C}) = 0$,
- (b) if $0 < (\mu_{\psi})^{c}(\mathbb{C}) < \infty$, then

$$\dim A^2_{e^{-\psi}}(\mathbb{C}) = \max \left\{ n \in \mathbb{N} : n < (\mu_{\psi})^c(\mathbb{C})/4\pi \right\}.$$

Recall that $A^2(\mathbb{P}^1; \mathcal{O}(k))$ is isomorphic to $A^2_{\phi_k}(\mathbb{C})$, and that Lemma 2.4 states that dim $A^2(\mathbb{P}^1, \mathcal{O}(k)) = \max\{0, k+1\}$. Using [2, Theorem 2.6], we can now rediscover Lemma 2.4 for $k \ge -2$. In particular, for $k \ge -2$, the function $\psi := -\ln(\phi_k)$ is subharmonic since

$$(\Delta \psi)(z) = 4(k+2)(1+|z|^2)^{-2}.$$

It follows that $(\mu_{\psi})^c(z) = 4(k+2)(1+|z|^2)^{-2}\lambda(z)$, so that $(\mu_{\psi})^c(\mathbb{C}) = 4\pi(k+2)$. By (a) and (b) above, it follows that dim $A_{\phi_k}^2(\mathbb{C}) = \max\{0, k+1\}$.

As in the case of unweighted Bergman spaces, the dimension of a weighted Bergman space for an open set in \mathbb{C} may be determined through the polarity of its complement as follows.

Corollary 3.3 Let $K \subset \mathbb{C}$ be a closed set, ψ a subharmonic function on \mathbb{C} . Then the following hold.

- If K is nonpolar, then A²_{e^{-ψ}} (ℂ \ K) is infinite dimensional.
 If K is polar, then A²_{e^{-ψ}} (ℂ \ K) is isomorphic to A²_{e^{-ψ}} (ℂ).

Proof To prove (1), we note first that Lemma 3.2 still holds if ψ_1 is subharmonic, not necessarily smooth, on Ω as long as the open set $U \in \Omega$ may be chosen such that $e^{-\psi_1}$ is integrable on U.

Next, we note that Propositions 2.1 and 2.2 in [10] imply that $e^{-\psi}$ is integrable except near finitely many points. Thus, we may choose an open set U in $\mathbb{C} \setminus K$ such that $e^{-\psi}$ is integrable on U. Now, choose ψ_2 equal to ψ_* as defined above (3.5). Then ψ_2 is a bounded above, strictly subharmonic, smooth function on $\mathbb{C} \setminus K$. Thus, with $\psi_1 := \psi$, the claim follows from the above mentioned variation of Lemma 3.2



For the proof of (2) we first note that ψ is upper semicontinuous, which implies that ψ is bounded from above on any compact subset of $\mathbb C$. Write Ω for $\mathbb C\setminus K$, and let $f\in A^2_{e^{-\psi}}(\Omega)$. It then follows that $f\in A^2(\Omega\cap\mathbb D(0,R))$ for any R>0. We now may proceed as in the paragraph following (3.1) and construct a function $F\in \mathcal O(\mathbb C)$ such that $F|_\Omega=f$. Since K is of Lebesgue measure 0, it then follows that $\|F\|_{L^2_{e^{-\psi}}(\mathbb C)}=\|f\|_{L^2_{e^{-\psi}}(\Omega)}$. Thus, $A^2_{e^{-\psi}}(\mathbb C\setminus K)$ is isomorphic to $A^2_{e^{-\psi}}(\mathbb C)$.

4 Finite Dimensional Bergman Spaces in \mathbb{P}^2

In this section, we show that the dichotomy displayed by Bergman spaces on \mathbb{P}^1 does not hold in higher dimensional projective spaces. In particular, for every holomorphic line bundle L on \mathbb{P}^2 , there exists a domain $D \subset \mathbb{P}^2$ such that the dimension of the Bergman space $A^2(D; L)$ is finite, but strictly larger than that of the space of global holomorphic sections of L. The examples in this section are entirely motivated by Wiegerinck's examples in [18].

In analogy with \mathbb{P}^1 , we use the coordinates $\zeta = [\zeta_0 : \zeta_1 : \zeta_2]$ on \mathbb{P}^2 . The open set $U^{(z,w)} = \mathbb{P}^2 \setminus \{\zeta : \zeta_0 = 0\}$ is identified with \mathbb{C}^2 via

$$\varphi_{(z,w)}: \zeta = [\zeta_0: \zeta_1: \zeta_2] \mapsto (z,w) = \left(\frac{\zeta_1}{\zeta_0}, \frac{\zeta_2}{\zeta_0}\right),$$

and (z, w) are referred to as the local coordinates on the affine chart $U^{(z,w)}$. Any holomorphic line bundle $\Pi: L \to \mathbb{P}^2$ is of the form $\mathcal{O}(k)$, for some $k \in \mathbb{Z}$, where $\mathcal{O}(k)$ is the line bundle associated to the divisor $k\{\ell\}$ for any fixed hyperplane $\ell \subset \mathbb{P}^1$. In the local coordinates (z, w), a global section of $\mathcal{O}(k)$ is given by a triplet, $s = (s_1, s_2, s_3)$, of holomorphic functions on \mathbb{C}^2 such that

$$s_1(z, w) = \begin{cases} z^k s_2(1/z, w/z), & \text{if } (z, w) \in \mathbb{C}^2 \setminus \{z = 0\}, \\ w^k s_3(z/w, 1/w), & \text{if } (z, w) \in \mathbb{C}^2 \setminus \{w = 0\}. \end{cases}$$
(4.1)

Similarly, a Hermitian metric on $\mathcal{O}(k)$ is given by a triplet, $h=(h_1,h_2,h_3)$, of smooth positive functions on \mathbb{C}^2 that satisfy compatibility conditions analogous to (4.1). In view of Lemma 2.2, we fix the following smooth volume form on \mathbb{P}^2 , and Hermitian metric on $\mathcal{O}(k)$, $k \in \mathbb{Z}$, respectively, in the local coordinates (z,w) on $U^{(z,w)}$:

$$\omega_{\text{FS}}(z, w) = (1 + |z|^2 + |w|^2)^{-3} dz \wedge d\overline{z} \wedge dw \wedge d\overline{w},$$

$$h_1(z, w) = (1 + |z|^2 + |w|^2)^{-k}.$$

Finally, recall the following description of the space of global sections of $\mathcal{O}(k)$, see [17, p. 81].



Lemma 4.1 Let $k \in \mathbb{Z}$. In the local coordinates (z, w) on the affine chart $U^{(z,w)}$,

$$A^{2}(\mathbb{P}^{2};\mathcal{O}(k)) = \Gamma(\mathbb{P}^{2};\mathcal{O}(k)) = \begin{cases} span\{z^{p}w^{q}: (p,q) \in \mathbb{N}^{2}, \ p+q \leq k\}, & if \ k \geq 0, \\ \{f \equiv 0\}, & if \ k < 0. \end{cases}$$

In particular, dim $A^2(\mathbb{P}^2; \mathcal{O}(k)) = \dim \Gamma(\mathbb{P}^2; \mathcal{O}(k)) = \max \left\{0, \frac{(k+1)(k+2)}{2}\right\}$.

Theorem 4.2 Given $k \in \mathbb{Z}$, there exists a domain $\Omega_k \subset \mathbb{P}^2$ such that

$$\dim A^2(\mathbb{P}^2; \mathcal{O}(k)) < \dim A^2(\Omega_k; \mathcal{O}(k)) < \infty.$$

In particular, dim $A^2(\Omega_k; \mathcal{O}(k)) = 1$ for k < -2, and dim $A^2(\Omega_k; \mathcal{O}(k)) = (k + 3)(k + 4)/2$ for $k \ge -2$.

Proof It suffices to produce a domain in the chart $U^{(z,w)}$, i.e., a domain $\Omega_k \subset \mathbb{C}^2$ such that dim $A^2_{\phi_k}(\mathbb{C}^2) < \dim A^2_{\phi_k}(\Omega_k) < \infty$, where

$$\phi_k(z, w) = (1 + |z|^2 + |w|^2)^{k+3}, \quad (z, w) \in \mathbb{C}^2.$$

The desired domain will be of the form $B \cup X_{\ell} \cup Y \cup Z_m$, for appropriately chosen $\ell, m \in \mathbb{N}$, where

$$\begin{split} B &= \left\{ (z,w) \in \mathbb{C}^2 : \max\{|z|,|w|\} < 2 \right\}, \\ X_{\ell} &= \left\{ (z,w) \in \mathbb{C}^2 : |z| > \sqrt{2}, \ |w| < 1/|z|^{\ell} \right\}, \quad \ell \in \mathbb{N}, \\ Y &= \left\{ (z,w) \in \mathbb{C}^2 : |w| > \sqrt{2}, \ |z| < 1/|w| \right\}, \\ Z_m &= \left\{ (z,w) \in \mathbb{C}^2 : |z|^2 + |w|^2 > 2, \ \left| |z| - |w| \right| < \frac{1}{(|z| + |w|)^m} \right\}, \quad m \in \mathbb{N} \setminus \{0,1\}. \end{split}$$

In the rest of the proof, $A_{k,\text{mon}}^2(\Omega)$ denotes the set of monomials in $A_{\phi_k}^2(\Omega)$, for any $\Omega \subset \mathbb{C}^2$. We claim that

$$A_{k,\text{mon}}^{2}(B) = \{ z^{p} w^{q} : (p,q) \in \mathbb{N}^{2} \}, \tag{4.2}$$

$$A_{k,\text{mon}}^2(X_\ell) = \{ z^p w^q : (p,q) \in \mathbb{N}^2, \ p - \ell q \le \ell + k + 1 \}, \tag{4.3}$$

$$A_{k \text{ mon}}^{2}(Y) = \{ z^{p} w^{q} : (p, q) \in \mathbb{N}^{2}, \ q - p \le k + 2 \}, \tag{4.4}$$

$$A_{k,\text{mon}}^2(Z_m) = \{ z^p w^q : (p,q) \in \mathbb{N}^2, \ 2(p+q) \le m + 2k + 2 \}.$$
 (4.5)

Assuming (4.2)–(4.5) for the moment, set

$$\Omega_k = \begin{cases} B \cup X_1 \cup Y \cup Z_2, & \text{if } k \ge -2, \\ B \cup X_{1-2(k+2)} \cup Y \cup Z_{-2(2k+3)}, & \text{if } k < -2. \end{cases}$$
(4.6)



Since Ω_k is Reinhardt for any $k \in \mathbb{Z}$, $A_{k,\text{mon}}^2(\Omega_k)$ is a basis for $A_{\phi_k}^2(\Omega_k)$. Thus, dim $A_{\phi_k}^2(\Omega_k)$ is the cardinality of $A_{k,\text{mon}}^2(\Omega_k)$. From (4.6), if $k \geq -2$,

$$A_{k,\text{mon}}^2(\Omega_k) = \{z^p w^q : (p,q) \in \mathbb{N}^2, |p-q| \le k+2, p+q \le k+2\},$$

which is of cardinality $\frac{(k+3)(k+4)}{2}$. On the other hand, if k < -2, i.e., k+2 = -t for some $t \ge 1$, then

$$A_{k,\text{mon}}^2(\Omega_k) = \left\{ z^p w^q : (p,q) \in \mathbb{N}^2, \ \frac{p-t}{1+2t} \le q \le p-t, \ p+q \le t \right\} = \{z^t\}.$$

Thus, recalling Lemma 4.1, $\dim A^2_{\phi_k}(\mathbb{C}^2) < \dim A^2_{\phi_k}(\Omega_k) < \infty$, in either case. It now remains to prove (4.2)–(4.5). We use the notation $L \approx M$ for $L, M \in \mathbb{R}$ to mean that there exist constants c, d > 0 such that $cM \le L \le dM$. Since (4.2) is clear, and (4.4) follows from (4.3), we only need to prove (4.3) and (4.5). To find the monomials which are contained in $A_{\phi_k}^2(X_\ell)$, consider

$$\widetilde{X}_{\ell} = \left\{ (r, s) : r > \sqrt{2}, 0 < s < 1/r^{\ell} \right\}.$$

For $p, q, \ell \in \mathbb{N}$, we get

$$||z^p w^q||_{A^2_{\phi_k}(X_\ell)}^2 = (2\pi)^2 \int_{\widetilde{X}_\ell} \frac{r^{2p+1} s^{2q+1}}{(1+r^2+s^2)^{3+k}} \; ds \; dr =: (2\pi)^2 \cdot \mathcal{J}.$$

To determine the range of p and q for which \mathcal{J} is finite for a given ℓ , we consider \widetilde{X}_{ℓ} in polar coordinates, i.e.,

$$\widetilde{X}_{\ell} = \{ (R\cos\theta, R\sin\theta) : R\cos\theta > \sqrt{2}, \ R^{\ell+1}\sin\theta(\cos\theta)^{\ell} < 1, \theta \in (0, \pi/2) \},$$

and introduce

$$\mathcal{J}(a,g) := \int_{a}^{\infty} \int_{0}^{g(R)} \frac{R^{2p+2q+2}}{(1+R^2)^{3+k}} (\cos \theta)^{2p+1} (\sin \theta)^{2q+1} R \ d\theta \ dR,$$

where a > 0 is a constant and g is a positive, continuous function. We notice that $\theta \in (0, \pi/4)$ on \widetilde{X}_{ℓ} . It then follows that

$$\mathcal{J}(2, \arcsin(R^{-\ell-1})) \le \mathcal{J} \le \mathcal{J}(\sqrt{2}, \arcsin(\sqrt{2}^{\ell} R^{-\ell-1})),$$

because, firstly, any point in \widetilde{X}_{ℓ} satisfies $R > \sqrt{2}$ and $\theta \in (0, \arcsin(\sqrt{2}^{\ell} R^{-\ell-1}))$, and, secondly, any point with R > 2 and $\theta \in (0, \arcsin(R^{-\ell-1}))$ is a point in \widetilde{X}_{ℓ} . Hence, it



remains to estimate $\mathcal{J}(a, \arcsin(cR^{-\ell-1})$ for $c \in \{1, \sqrt{2}^{\ell}\}$. Again, since $\theta \in (0, \pi/4)$ on \widetilde{X}_{ℓ} , it follows that

$$\mathcal{J}(a, \arcsin(cR^{-\ell-1})) = \int_{a}^{\infty} \frac{R^{2p+2q+2}}{(1+R^2)^{3+k}} \times \left(\int_{0}^{\arcsin(cR^{-\ell-1})} (\cos\theta)^{2p+1} (\sin\theta)^{2q+1} d\theta \right) R dR$$

$$\approx \int_{a}^{\infty} \frac{R^{2p+2q+2}}{(1+R^2)^{3+k}} \times \left(\int_{0}^{\arcsin(cR^{-\ell-1})} \cos\theta (\sin\theta)^{2q+1} d\theta \right) R dR$$

$$\approx \int_{a}^{\infty} R^{2p+2q+3-(\ell+1)(2q+2)-2(3+k)} dR,$$

which converges if and only if $2p + 2q + 3 - (\ell + 1)(2q + 2) - 2(3 + k) < -1$. It follows that $z^p w^q \in A^2_{\phi_k}(X_\ell)$ if and only if

$$p - \ell q \le k + \ell + 1. \tag{4.7}$$

In a similar fashion, we find the monomials in $A_{\phi_k}^2(Z_m)$. For that, we write

$$\widetilde{Z}_m = \left\{ (r, s) : r^2 + s^2 > 2, |r - s|(r + s)^m < 1 \right\}$$

and note that for $p, q, m \in \mathbb{N}$

$$||z^p w^q||_{A_{\phi_k}^2(Z_m)}^2 = (2\pi)^2 \int_{\widetilde{Z}_m} \frac{r^{2p+1} s^{2q+1}}{(1+r^2+s^2)^{3+k}} dr ds := (2\pi)^2 \cdot \mathcal{I}.$$

As before, we introduce polar coordinates on \widetilde{Z}_m . It is straightforward to check that then

$$\begin{split} \widetilde{Z}_m &= \left\{ (R\cos\theta, R\sin\theta) : R > \sqrt{2}, \; \left| \sin\left(\frac{\pi}{4} - \theta\right) \right| \cos^m\left(\frac{\pi}{4} - \theta\right) < (\sqrt{2}R)^{-m-1}, \theta \in (0, \frac{\pi}{2}) \right\} \\ &= \left\{ \left(R\cos\left(\frac{\pi}{4} - \psi\right), R\sin\left(\frac{\pi}{4} - \psi\right)\right) : R > \sqrt{2}, \; \left| \sin\psi \right| \cos^m\psi < (\sqrt{2}R)^{-m-1}, \psi \in (-\frac{\pi}{4}, \frac{\pi}{4}) \right\}. \end{split}$$

We work in the coordinates (R, ψ) to estimate \mathcal{I} , utilizing integrals of the form

$$\mathcal{I}(g) := \int_{\sqrt{2}}^{\infty} \int_{-g(R)}^{g(R)} \frac{R^{2p+2q+3}}{(1+R^2)^{3+k}} \cos^{2p+1} \left(\frac{\pi}{4} - \psi\right) \sin^{2q+1} \left(\frac{\pi}{4} - \psi\right) d\psi dR$$

for some positive, continuous function g. Since $\psi \in [-\pi/4, \pi/4)$, it follows that

$$\mathcal{I}\left(\arcsin\left((\sqrt{2}R)^{-m-1}\right)\right) \leq \mathcal{I} \leq \mathcal{I}\left(\arcsin\left(\sqrt{2}R^{-m-1}\right)\right).$$



Since $m \ge 2$, both trigonometric functions in the definition of $\mathcal{I}(g)$ are strictly positive for $|\psi| \le |\arcsin(\sqrt{2}R^{-m-1})|$. Therefore,

$$\mathcal{I}\left(\arcsin(cR^{-m-1})\right) \approx \int_{\sqrt{2}}^{\infty} \int_{-\arcsin(cR^{-m-1})}^{\arcsin(cR^{-m-1})} \frac{R^{2p+2q+3}}{(1+R^2)^{3+k}} \ d\psi \ dR$$

for $c \in {\sqrt{2}, \sqrt{2}^{-m-1}}$. Moreover, $|t/2| \le |\sin t| \le |t|$ holds for $|t| \le \pi/4$. Thus

$$\mathcal{I}\left(\arcsin(cR^{-m-1})\approx \int_{\sqrt{2}}^{\infty}R^{2p+2q+3-(m+1)-2(3+k)}dR,\right.$$

which converges if and only if

$$2(p+q) \le m + 2k + 2. \tag{4.8}$$

This completes the proof of (4.5), and thus, of Theorem 4.2.

Acknowledgements The authors are grateful to Peter Pflug for his helpful comments on Lemma 3.1.

Declarations

Conflict of interest The authors have no conflict of interest to declare that are relevant to the content of this article.

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