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Solving the Kohn Laplacian on asymptotically flat CR manifolds of dimension 3



MATHEMATICS

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Chin-Yu Hsiao^{a,*,1}, Po-Lam Yung^{b,2}

 ^a Institute of Mathematics, Academia Sinica, 6F, Astronomy-Mathematics Building, No. 1, Sec. 4, Roosevelt Road, Taipei 10617, Taiwan
 ^b Department of Mathematics, the Chinese University of Hong Kong, Shatin, Hong Kong

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ABSTRACT

Let $(\hat{X}, T^{1,0}\hat{X})$ be a compact orientable CR embeddable three dimensional strongly pseudoconvex CR manifold, where $T^{1,0}\hat{X}$ is a CR structure on \hat{X} . Fix a point $p \in \hat{X}$ and take a global contact form $\hat{\theta}$ so that $\hat{\theta}$ is asymptotically flat near p. Then $(\hat{X}, T^{1,0}\hat{X}, \hat{\theta})$ is a pseudohermitian 3-manifold. Let $G_p \in C^{\infty}(\hat{X} \setminus \{p\}), G_p > 0$, with $G_p(x) \sim \vartheta(x, p)^{-2}$ near p, where $\vartheta(x, y)$ denotes the natural pseudohermitian distance on \hat{X} . Consider the new pseudohermitian 3-manifold with a blow-up of contact form $(\hat{X} \setminus \{p\}, T^{1,0}\hat{X}, G_p^2\hat{\theta})$ and let \Box_b denote the corresponding Kohn Laplacian on $\hat{X} \setminus \{p\}$.

In this paper, we prove that the weighted Kohn Laplacian $G_p^2 \square_b$ has closed range in L^2 with respect to the weighted volume form $G_p^2 \hat{\theta} \wedge d\hat{\theta}$, and that the associated partial inverse and the Szegö projection enjoy some regularity properties near p. As an application, we prove the existence of some special functions in the kernel of \square_b that grow at a specific rate at p. The existence of such functions provides an important ingredient for the proof of a positive mass theorem

* Corresponding author.

E-mail addresses: chsiao@math.sinica.edu.tw (C.-Y. Hsiao), plyung@math.cuhk.edu.hk (P.-L. Yung).

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in 3-dimensional CR geometry by Cheng, Malchiodi and Yang [5]. © 2015 Elsevier Inc. All rights reserved.

1. Introduction

1.1. Motivation from CR geometry

The study described in this paper was motivated by that on a positive mass theorem in 3-dimensional CR geometry by Cheng, Malchiodi and Yang [5], where one needs to find some special functions in the kernel of the Kohn Laplacian that grow at a specific rate at a given point on an asymptotically flat pseudohermitian 3-manifold. We begin by giving a brief description of the relevance of our result with their work below.

Consider a compact orientable 3-dimensional strongly pseudoconvex CR manifold $(\hat{X}, T^{1,0}\hat{X})$ with CR structure $T^{1,0}\hat{X}$. We assume throughout that it is CR embeddable in some \mathbb{C}^N . By choosing a contact form $\hat{\theta}_0$ on \hat{X} that is compatible with its CR structure, one can make $(\hat{X}, T^{1,0}\hat{X}, \hat{\theta}_0)$ a pseudohermitian 3-manifold; in particular, one can define a Hermitian inner product on $T^{1,0}\hat{X}$, by

$$\langle Z_1 | Z_2 \rangle_{\hat{\theta}_0} = \frac{1}{2} d\hat{\theta}_0(Z_1, i\overline{Z}_2).$$

Now fix $p \in \hat{X}$. By conformally changing the contact form, we may find another contact form $\hat{\theta}$ (which is a multiple of $\hat{\theta}_0$ by a positive smooth function), so that near p, there exists CR normal coordinates (z, t). In other words, the contact form $\hat{\theta}$ and the coordinates (z, t) are chosen, so that

- (i) the point p corresponds to (z,t) = (0,0);
- (ii) one can find a local section \hat{Z}_1 of $T^{1,0}\hat{X}$ near p, with $\langle \hat{Z}_1 | \hat{Z}_1 \rangle_{\hat{\theta}} = 1$, such that \hat{Z}_1 admits an expansion near p as described in (1.31) below; and
- (iii) the Reeb vector field \hat{T} with respect to $\hat{\theta}$ admits an expansion as described in (1.32).

Then $(\hat{X}, T^{1,0}\hat{X}, \hat{\theta})$ is another pseudohermitian 3-manifold. Assume that this pseudohermitian 3-manifold is of positive Tanaka–Webster class: this means that the lowest eigenvalue of the conformal sublaplacian

$$L_b := -4\Delta_b + R$$

is strictly positive. Here $R = R_{\hat{\theta}}$ is the Tanaka–Webster curvature of \hat{X} , and Δ_b is the sublaplacian on \hat{X} . (The above assumption on L_b will hold when e.g. $R_{\hat{\theta}}$ is strictly positive on \hat{X} .) Then L_b is invertible, so one can write down the Green's function G_p of L_b with pole at p. We normalize G_p so that

$$L_b G_p = 16\delta_p.$$

Then near p, G_p admits the following expansion

$$G_p = \frac{1}{2\pi}\hat{\rho}^{-2} + A + f, \quad f \in \mathcal{E}(\hat{\rho}^1),$$
(1.1)

where A is some real constant, $\hat{\rho}^4(z,t) = |z|^4 + t^2$, and for $m \in \mathbb{R}$, $\mathcal{E}(\hat{\rho}^m)$ denotes, roughly speaking, the set of all smooth functions $g \in C^{\infty}(\hat{X} \setminus \{p\})$ such that $|g(z,t)| \leq \hat{\rho}(z,t)^{m-p-q-2r}$ near p, along with some suitable control of the growth of derivatives near p (see (1.34) for the precise meaning of the Fréchet space $\mathcal{E}(\hat{\rho}^m)$).

Let now

$$X = \hat{X} \setminus \{p\}$$
 and $\theta = G_p^2 \hat{\theta}$

and let $T^{1,0}X$ be the restriction of $T^{1,0}\hat{X}$ to X. Then $(X, T^{1,0}X, \theta)$ is a new non-compact pseudohermitian manifold, which we think of as the blow-up of our original $(\hat{X}, T^{1,0}\hat{X}, \hat{\theta})$ at p. We say that this pseudohermitian manifold is asymptotically flat, since under an inversion of coordinates, X has asymptotically the geometry of the Heisenberg group at infinity. We note that the Tanaka–Webster scalar curvature R_{θ} of X is identically zero, since the conformal factor G_p we used is the Green's function for the conformal sublaplacian on \hat{X} . Let $\bar{\partial}_b$ and \Box_b denote the associated tangential Cauchy–Riemann operator and Kohn Laplacian respectively. In [5], Cheng, Malchiodi and Yang introduced the pseudohermitian p-mass for $(X, T^{1,0}X, \theta)$, given by

$$m(\theta) := \lim_{\Lambda \to 0} i \int_{\{\hat{\rho} = \Lambda\}} \omega_1^1 \wedge \theta,$$

where ω_1^1 stands for the connection form of the given pseudohermitian structure. Then they proved that there is a specific $\tilde{\beta} \in \mathcal{E}(\hat{\rho}^{-1})$, with $\Box_b \tilde{\beta} = \mathcal{E}(\hat{\rho}^4)$, such that

$$m(\theta) = -\frac{3}{2} \int_{X} \left| \Box_{b} \widetilde{\beta} \right|^{2} \theta \wedge d\theta + 3 \int_{X} |\widetilde{\beta}_{,\bar{1}\bar{1}}|^{2} \theta \wedge d\theta + \frac{3}{4} \int_{X} \overline{\widetilde{\beta}} \cdot P \widetilde{\beta} \theta \wedge d\theta, \qquad (1.2)$$

where $\tilde{\beta}_{,\overline{11}}$ is some derivative of the function $\tilde{\beta}$, and P is the CR Paneitz operator of $(X, T^{1,0}X, \theta)$. (Note $R_{\theta} = 0$ in our current set-up, so the term involving R_{θ} in the corresponding identity of mass in [5] is not present above.) Moreover, it was shown that (1.2) holds for any β in place of $\tilde{\beta}$, as long as $\tilde{\beta} - \beta \in \mathcal{E}(\hat{\rho}^{1+\delta})$, and $\Box_b \beta = \mathcal{E}(\hat{\rho}^{3+\delta})$ for some $\delta > 0$. Thus, if we could find such a β in the *kernel* of \Box_b , then under the assumption that the CR Paneitz operator P is non-negative, one can conclude that the mass $m(\theta)$ is non-negative. The construction of such β is the motivation of the current paper. (See Corollary 1.2 in the next subsection.)

Classically, if one wants to solve \Box_b on say a compact CR manifold, one proceeds by showing first that \Box_b extends to a closed linear operator on L^2 , and that this extended

 \Box_b has closed range on L^2 . Then one solves \Box_b in a weak sense, and shows that the solution is classical if the right hand side of the equation is smooth. This strategy does not directly apply in our situation, since our CR manifold X is non-compact. In fact, the natural volume form on X is given by $\theta \wedge d\theta$, and even if we extend \Box_b so that it becomes a closed linear operator on $L^2(\theta \wedge d\theta)$, this operator may not have closed range in $L^2(\theta \wedge d\theta)$, as is seen in some simple examples (e.g. when \hat{X} is the unit sphere in \mathbb{C}^2).

To overcome this difficulty, we introduce in this paper a weighted volume form

$$m_1 := G_p^{-2}\theta \wedge d\theta,$$

as well as a weighted Kohn Laplacian, namely

$$\Box_{b,1} := G_p^2 \Box_b.$$

We will show that $\Box_{b,1}$ extends as a densely defined closed linear operator

$$\square_{b,1}: \text{Dom} \square_{b,1} \subset L^2(m_1) \to L^2(m_1),$$

and this extended operator has closed range in $L^2(m_1)$. (See Theorem 1.3 below.) Here $L^2(m_1)$ is the space of L^2 functions with respect to the volume form m_1 . As a result, we have the following L^2 decomposition:

$$\square_{b,1}N + \Pi = I$$
 on $L^2(m_1)$

Here $N : L^2(m_1) \to \text{Dom} \square_{b,1}$ is the partial inverse of $\square_{b,1}$ and $\Pi : L^2(m_1) \to (\text{Ran} \square_{b,1})^{\perp}$ is the orthogonal projection onto the orthogonal complement of the range of $\square_{b,1}$ in $L^2(m_1)$. We will show that for every $0 < \delta < 2$, N and Π can be extended continuously to

$$N : \mathcal{E}(\hat{\rho}^{-2+\delta}) \to \mathcal{E}(\hat{\rho}^{\delta}),$$

$$\Pi : \mathcal{E}(\hat{\rho}^{-2+\delta}) \to \mathcal{E}(\hat{\rho}^{-2+\delta})$$
(1.3)

(see Theorem 1.4 below). Hence

$$\Box_{b,1}N + \Pi = I \quad \text{on} \quad \mathcal{E}(\hat{\rho}^{-2+\delta}), \tag{1.4}$$

for every $0 < \delta < 2$. Now, let $\tilde{\beta}$ be as in (1.2). Put

$$f := \Box_{b,1}\widetilde{\beta} = G_p^2 \Box_b \widetilde{\beta}$$

Then by the expansion of G_p and the assumption on $\Box_b \tilde{\beta}$, we have $f \in \mathcal{E}(\hat{\rho}^{-1+\delta})$, for every $0 < \delta < 1$. From (1.3), we know that Πf is well-defined and $Nf \in \mathcal{E}(\hat{\rho}^{1+\delta})$, for every $0 < \delta < 1$. Moreover, we will show, in Theorem 1.5 below, that

$$\Pi f = 0$$

Thus from (1.4), we have

$$\Box_{b,1}(Nf) = f = \Box_{b,1}\beta.$$

If we put

$$\beta := \beta - Nf,$$

then

$$\Box_b \beta = 0 \text{ and } \beta - \widetilde{\beta} \in \mathcal{E}(\widehat{\rho}^{1+\delta}),$$

for every $0 < \delta < 1$. With this we have achieved our goal.

It turns out that a large part of our analysis does not depend on the fact that G_p is the Green's function of a conformal sublaplacian. All that we need is that G_p admits an expansion as in (1.1), that it is smooth on X, and that it is positive everywhere on X. We will formulate our result in this framework in the next subsection.

We expect that it is possible to approach the same problem by proceeding via L^p spaces rather than weighted L^2 spaces. In [13], we solved the \Box_b equation in some L^p spaces in a special case.

The operator $\Box_{b,1}$ can be seen as the Kohn Laplacian on the non-compact CR manifold $X = \hat{X} \setminus \{p\}$ defined with respect to the natural CR structure $T^{1,0}\hat{X}$ and the "singular" volume form m_1 . The coefficients of $\Box_{b,1}$ are smooth on X but singular at p. This work can be seen as a first study of this kind of "singular Kohn Laplacians". It will be quite interesting to develop some kind of "singular" functional calculus for pseudodifferential operators and Fourier integral operators and establish a completely microlocal analysis for $\Box_{b,1}$ along the lines of Beals and Greiner [1], Boutet de Monvel and Sjöstrand [3] and [12]. We hope that the "singular Kohn Laplacians" will be interesting for analysts.

1.2. Our main result

Let us now formulate our main results in their full generality. Consider a compact orientable 3-dimensional strongly pseudoconvex pseudohermitian manifold \hat{X} , with CR structure $T^{1,0}\hat{X}$ and contact form $\hat{\theta}_0$. We assume throughout that it is CR embeddable in some \mathbb{C}^N . By conformally changing the contact form $\hat{\theta}_0$, we may find another contact form $\hat{\theta}$, so that near p, one can find CR normal coordinates (z, t) as described in the previous subsection. We will write

$$\hat{\rho}(z,t) = (|z|^4 + t^2)^{1/4},$$

and for every $m \in \mathbb{R}$, we can define a Fréchet function space $\mathcal{E}(\hat{\rho}^m)$ as in (1.34).

Now fix a point $p \in \hat{X}$, and let

$$X = \hat{X} \setminus \{p\}.$$

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We fix from now on an everywhere positive function $G_p \in C^{\infty}(X)$, such that G_p admits an expansion

$$G_p = \frac{1}{2\pi}\hat{\rho}^{-2} + A + f, \quad f \in \mathcal{E}(\hat{\rho}^1), \tag{1.5}$$

where A is some real constant. (Again G_p need not be the Green's function of the conformal sublaplacian any more.) Define now

$$\theta = G_p^2 \hat{\theta}.$$

Then $(X, T^{1,0}\hat{X}, \theta)$ is a non-compact pseudohermitian 3-manifold, which we think of as the blow-up of the original \hat{X} . The θ defines for us a pointwise Hermitian inner product on $T^{1,0}X$, given by

$$\langle Z_1 | Z_2 \rangle_{\theta} = \frac{1}{2} d\theta (Z_1, i \overline{Z}_2);$$

we denote the dual pointwise inner product on the space (0, 1) forms on X by the same notation $\langle \cdot | \cdot \rangle_{\theta}$. Let $\overline{\partial}_b$ be the tangential Cauchy–Riemann operator on X. This is defined depending only on the CR structure on X. Now there is a natural volume form on X, given by

$$m := \theta \wedge d\theta.$$

This induces an inner product on functions on X, given by

$$(f|g)_m = \int\limits_X f\bar{g} \, m,$$

and an inner product on (0, 1) forms on X, given by

$$(\alpha|\beta)_{m,\theta} = \int\limits_X \langle \alpha|\beta \rangle_{\theta} \, m.$$

We write $\bar{\partial}_b^{*,f}$ for the formal adjoint of $\bar{\partial}_b$ under these two inner products. In other words, $\bar{\partial}_b^{*,f}$ satisfies

$$(\overline{\partial}_b u | v)_{m,\theta} = (u | \overline{\partial}_b^{*,f} v)_m$$

for all functions u and (0,1) forms v on X that are smooth with compact support. We can now define the Kohn Laplacian on X, namely

$$\Box_b := \overline{\partial}_b^{*,f} \overline{\partial}_b,$$

at least on smooth functions with compact support on X; then

$$(\Box_b u \mid f)_m = (u \mid \Box_b f)_m$$

for all functions u, f on X that are smooth with compact support, so we can extend \Box_b to distributions on X by duality.

Our goal is then to solve a specific equation involving \Box_b . First, let $\chi(z,t)$ be a smooth function with compact support on \hat{X} , so that its support is contained in the local coordinate chart given by the CR normal coordinates (z,t), and that it is identically 1 in a neighborhood of p. Let

$$\beta_0 = \chi(z,t) \frac{i\overline{z}}{|z|^2 - it} \in \mathcal{E}(\hat{\rho}^{-1}).$$

Then $\Box_b \beta_0 \in \mathcal{E}(\hat{\rho}^3)$. Furthermore, as was shown in [5], there exists $\beta_1 \in \mathcal{E}(\hat{\rho}^1)$, such that if

$$\tilde{\beta} := \beta_0 + \beta_1,$$

then

$$F := \Box_b \tilde{\beta} \in \mathcal{E}(\hat{\rho}^4).$$

Actually, in what follows, all we will use is that $F \in \mathcal{E}(\hat{\rho}^{3+\delta})$ for all $0 < \delta < 1$. Our main theorem can now be stated as follows:

Theorem 1.1. Let F be as defined above. Then there exists a smooth function u on X, such that $u \in \mathcal{E}(\hat{\rho}^{1+\delta})$ for any $0 < \delta < 1$, and

$$\Box_b u = F.$$

By taking $\beta := \tilde{\beta} - u$, we then have:

Corollary 1.2. There exists $\beta \in \mathcal{E}(\hat{\rho}^{-1})$ with $\beta - \tilde{\beta} \in \mathcal{E}(\hat{\rho}^{1+\delta})$ for any $0 < \delta < 1$, such that

$$\Box_b \beta = 0.$$

This provides a key tool in the proof of a positive mass theorem in 3-dimensional CR geometry in the work of Cheng, Malchiodi and Yang [5], as was explained in the last subsection.

Some remarks are in order. The first is about numerology. Considerations of homogeneity shows that \Box_b takes a function in $\mathcal{E}(\hat{\rho}^k)$ to $\mathcal{E}(\hat{\rho}^{k+2})$. Thus the homogeneity above works out right; the only small surprise is that while β_0 is in $\mathcal{E}(\hat{\rho}^{-1})$, $\Box_b\beta_0$ is in $\mathcal{E}(\hat{\rho}^3)$, which is 2 orders better than expected. But that is a reflection of the fact that our β_0 has been chosen such that $\bar{\partial}_b \beta_0$ is almost annihilated by $\bar{\partial}_b^{*,f}$.

Next, β_1 above is an explicit correction in $\mathcal{E}(\hat{\rho}^1)$ such that $\Box_b\beta_1 \in \mathcal{E}(\hat{\rho}^3)$ cancels out the main contribution of $\Box_b\beta_0 \in \mathcal{E}(\hat{\rho}^3)$. This ensures that $F \in \mathcal{E}(\hat{\rho}^{3+\delta})$ for all $0 < \delta < 1$, which in turn guarantees that the β we construct in Corollary 1.2 is determined explicitly up to $\mathcal{E}(\hat{\rho}^{1+\delta})$ for all (in particular, for some) $0 < \delta < 1$. The latter is important in the proof of the positive mass theorem in [5], since any term in the expansion of β that is in $\mathcal{E}(\hat{\rho}^1)$ would enter into the calculation of mass in (1.2). But for the purpose of solving \Box_b in the current paper, the correction term β_1 is not essential; in particular, if $F_0 := \Box_b\beta_0 \in \mathcal{E}(\hat{\rho}^3)$, then our proof below carries over, and shows that there exists $u_0 \in \mathcal{E}(\hat{\rho}^1)$ such that $\Box_b u_0 = F_0$.

Finally, in Corollary 1.2, note that we do not claim $\overline{\partial}_b \beta = 0$. It is only $\Box_b \beta$ that vanishes, as can be shown by say the example when \hat{X} is the standard CR sphere in \mathbb{C}^2 .

1.3. Our strategy

As we mentioned earlier, the difficulty in establishing the above theorem is that the CR manifold we are working on, namely X, is non-compact; also, the natural measure on X, namely $m = \theta \wedge d\theta$, has infinite volume on X. Let $L^2(m)$ be the space of L^2 functions on X with respect to m. Even if we extend \Box_b to be a closed linear operator on $L^2(m) \to L^2(m)$, in general the extended \Box_b may not have closed range in $L^2(m)$. Thus the classical methods of solving \Box_b fail in our situation.

We thus proceed by introducing a weighted L^2 space, and a weighted Kohn Laplacian. Let

$$m_1 := G_p^{-2}\theta \wedge d\theta.$$

We define $L^2(m_1)$ to be the space of L^2 functions on X with respect to the inner product

$$(f|g)_{m_1} := \int\limits_X f\overline{g} \, m_1,$$

and define $L^2_{(0,1)}(m_1, \hat{\theta})$ to be the space of $L^2(0, 1)$ forms on X with respect to the inner product

$$(\alpha|\beta)_{m_1,\hat{\theta}} := \int\limits_X \langle \alpha|\beta \rangle_{\hat{\theta}} \, m_1$$

We extend the tangential Cauchy–Riemann operator so that

 $\text{Dom}\overline{\partial}_{b,1} := \{ u \in L^2(m_1) : \text{the distributional } \overline{\partial}_b \text{ of } u \text{ on } X \text{ is in } L^2_{(0,1)}(m_1, \hat{\theta}) \},\$

and define

$$\overline{\partial}_{b,1} u := \text{the distributional } \overline{\partial}_b \text{ of } u \text{ on } X$$

if $u \in \text{Dom}\overline{\partial}_{b,1}$. Then

$$\overline{\partial}_{b,1}$$
: Dom $\overline{\partial}_{b,1} \subset L^2(m_1) \to L^2_{(0,1)}(m_1,\hat{\theta}),$

is a densely defined closed linear operator. Let

$$\overline{\partial}_{b,1}^*$$
: Dom $\overline{\partial}_{b,1}^* \subset L^2_{(0,1)}(m_1,\hat{\theta}) \to L^2(m_1)$

be its adjoint. Let $\square_{b,1}$ denote the Gaffney extension of the singular Kohn Laplacian given by

$$\operatorname{Dom} \Box_{b,1} = \left\{ s \in L^2(m_1) : s \in \operatorname{Dom} \overline{\partial}_{b,1}, \ \overline{\partial}_{b,1} s \in \operatorname{Dom} \overline{\partial}_{b,1}^* \right\} \,,$$

and $\Box_{b,1}s = \overline{\partial}_{b,1}^*\overline{\partial}_{b,1}s$ for $s \in \text{Dom} \Box_{b,1}$. By a result of Gaffney, $\Box_{b,1}$ is a non-negative self-adjoint operator (see [19, Prop. 3.1.2]). We extend $\Box_{b,1}$ to distributions on X by

$$\Box_{b,1} : \mathscr{D}'(X) \to \mathscr{D}'(X),$$
$$(\Box_{b,1}u \mid f)_{m_1} = (u \mid \Box_{b,1}f)_{m_1}, \ u \in \mathscr{D}'(X), \ f \in C_0^\infty(X).$$

This is well-defined, since if f is a test function on X, then so is $\Box_{b,1}f$. One can show

$$\Box_{b,1} u = G_p^2 \Box_b u, \quad \forall u \in \mathscr{D}'(X).$$

In fact, by duality, it suffices to check this for a test function $u \in C_0^{\infty}(X)$. If u is as such, then for any test function $v \in C_0^{\infty}(X)$,

$$(\Box_{b,1}u \mid v)_{m_1} = (\overline{\partial}_{b,1}u \mid \overline{\partial}_{b,1}v)_{m_1,\hat{\theta}}$$
$$= (\overline{\partial}_b u \mid \overline{\partial}_b v)_{m_1,\hat{\theta}}$$
$$= (\overline{\partial}_b u \mid G_p^{-2}\overline{\partial}_b v)_{m,\hat{\theta}}$$
$$= (\overline{\partial}_b u \mid \overline{\partial}_b v)_{m,\theta}$$
$$= (\Box_b u \mid v)_m$$
$$= (G_p^2 \Box_b u \mid v)_{m_1},$$

and the desired equality follows.

Thus solving \Box_b is essentially the same as solving for $\Box_{b,1}$, and it is the latter that forms the heart of our paper.

The key here is then three-fold, as is represented by the next three theorems. First we will show that

Theorem 1.3. $\Box_{b,1}$: Dom $\Box_{b,1} \subset L^2(m_1) \to L^2(m_1)$ has closed range in $L^2(m_1)$.

Once this is shown, we have the following L^2 decomposition:

$$\Box_{b,1}N + \Pi = I \quad \text{on} \quad L^2(m_1).$$

Here $N : L^2(m_1) \to \text{Dom} \square_{b,1}$ is the partial inverse of $\square_{b,1}$ and $\Pi : L^2(m_1) \to (\text{Ran} \square_{b,1})^{\perp}$ is the orthogonal projection onto $(\text{Ran} \square_{b,1})^{\perp}$. We now need:

Theorem 1.4. For every $0 < \delta < 2$, Π and N can be extended continuously to

$$\Pi : \mathcal{E}(\hat{\rho}^{-2+\delta}) \to \mathcal{E}(\hat{\rho}^{-2+\delta})$$
$$N : \mathcal{E}(\hat{\rho}^{-2+\delta}) \to \mathcal{E}(\hat{\rho}^{\delta}).$$

It then follows that

$$\Box_{b,1}N + \Pi = I \text{ on } \mathcal{E}(\hat{\rho}^{-2+\delta}),$$

for every $0 < \delta < 2$. Now, let F be as in Theorem 1.1. Put

$$f := G_n^2 F$$

Then $f \in \mathcal{E}(\hat{\rho}^{-1+\delta})$, for every $0 < \delta < 1$. From Theorem 1.4, we know that Πf is well-defined; we will show that

Theorem 1.5.

 $\Pi f = 0.$

It then follows that u := Nf satisfies

$$u \in \mathcal{E}(\hat{\rho}^{1+\delta}), \quad \Box_{b,1}u = f = G_p^2 F.$$

From the relation between $\Box_{b,1}$ and \Box_b , we obtain the desired conclusion in Theorem 1.1.

1.4. Outline of proofs

To prove the theorems in the previous subsection, we need to introduce two other Kohn Laplacians, which we denote by $\hat{\Box}_b$ and $\widetilde{\Box}_b$, as follows.

First, $\hat{\Box}_b$ is the natural Kohn Laplacian on $(\hat{X}, T^{1,0}\hat{X}, \hat{\theta})$. There we have the natural measure

$$\hat{m} := \hat{\theta} \wedge d\hat{\theta}.$$

One can then define $L^2(\hat{m})$ to be the space of L^2 functions on \hat{X} with respect to the inner product

$$(f|g)_{\hat{m}} := \int\limits_{\hat{X}} f \bar{g} \, \hat{m},$$

and define $L^2_{(0,1)}(\hat{m},\hat{\theta})$ to be the space of L^2 (0,1) forms on \hat{X} with respect to the inner product

$$(\alpha|\beta)_{\hat{m},\hat{\theta}} := \int_{\hat{X}} \langle \alpha|\beta \rangle_{\hat{\theta}} \, \hat{m}.$$

We extend the tangential Cauchy–Riemann operator so that

$$\text{Dom}\overline{\bar{\partial}}_b := \{ u \in L^2(\hat{m}) : \text{the distributional } \overline{\partial}_b \text{ of } u \text{ on } \hat{X} \text{ is in } L^2_{(0,1)}(\hat{m}, \hat{\theta}) \}$$

and define

$$\overline{\partial}_b u :=$$
the distributional $\overline{\partial}_b$ of u on \hat{X}

if $u \in \text{Dom}\overline{\partial}_b$. Then

$$\hat{\overline{\partial}}_b$$
: Dom $\hat{\overline{\partial}}_b \subset L^2(\hat{m}) \to L^2_{(0,1)}(\hat{m},\hat{\theta})$

is a densely defined closed linear operator. Let

$$\hat{\overline{\partial}}_b^*: \operatorname{Dom} \hat{\overline{\partial}}_b^* \subset L^2_{(0,1)}(\hat{m}, \hat{\theta}) \to L^2(\hat{m})$$

be its adjoint. Let $\hat{\Box}_b$ denote the Gaffney extension of the Kohn Laplacian given by

$$\operatorname{Dom} \hat{\Box}_b = \{ s \in L^2(\hat{m}) : s \in \operatorname{Dom} \hat{\overline{\partial}}_b, \ \hat{\overline{\partial}}_b s \in \operatorname{Dom} \hat{\overline{\partial}}_b^* \}, \\ \hat{\Box}_b s = \hat{\overline{\partial}}_b^* \hat{\overline{\partial}}_b s \quad \text{for } s \in \operatorname{Dom} \hat{\Box}_b.$$

It is then a non-negative self-adjoint operator on $L^2(\hat{m})$ (see e.g. [19, Prop. 3.1.2]). The analysis of this $\hat{\Box}_b$ is very well-understood; see work of Kohn [16,17], Boas–Shaw [2], Christ [6,7] and Fefferman–Kohn [9] in the CR embeddable case, and work of Kohn– Rossi [18], Folland–Stein [10], Rothschild–Stein [26], Greiner–Stein [11], Nagel–Stein [22], Fefferman [8], Boutet de Monvel–Sjöstrand [3], Nagel–Stein–Wainger [23], Nagel–Rosay– Stein–Wainger [24,25] and Machedon [20,21] for some earlier work or related results. On the other hand, it is not very straightforward to reduce the analysis of $\Box_{b,1}$ to the analysis of this $\hat{\Box}_b$; we go through an intermediate Kohn Laplacian, which we denote by $\tilde{\Box}_b$. To introduce $\tilde{\Box}_b$, first we need to construct a special CR function ψ on \hat{X} , such that $\bar{\partial}_b \psi = 0$ on \hat{X} , $\psi \neq 0$ on X, and near p, we have

$$\psi(z,t) = 2\pi(|z|^2 + it) + \text{error},$$

where the error vanishes like $\hat{\rho}^4$ near p. (The precise construction is given in Section 4.) One can then define the following volume form on \hat{X} :

$$\tilde{m} := G_p^2 |\psi|^2 \hat{\theta} \wedge d\hat{\theta}.$$

Note that \hat{X} has finite volume with respect to this volume form, since $G_p^2 |\psi|^2$ is bounded near p. However, this volume form does not have a smooth density against \hat{m} ; this is one of the biggest sources of difficulties in what we do below. The key turns out to be the following: the asymptotics (1.5) we assumed of G_p allows us to obtain some crucial asymptotics of the density of \tilde{m} against \hat{m} near p. This in turn implies a crucial relation between the \Box_b we will introduce, and the \Box_b we defined above (see (1.9) below).

Now let $L^2(\widetilde{m})$ be the space of L^2 functions on \hat{X} with respect to the inner product

$$(f|g)_{\tilde{m}} := \int_{\hat{X}} f \bar{g} \, \tilde{m},$$

and define $L^2_{(0,1)}(\tilde{m},\hat{\theta})$ to be the space of $L^2(0,1)$ forms on \hat{X} with respect to the inner product

$$(\alpha|\beta)_{\tilde{m},\hat{\theta}} := \int_{\hat{X}} \langle \alpha|\beta \rangle_{\hat{\theta}} \, \tilde{m}.$$

We extend the tangential Cauchy-Riemann operator so that

$$\operatorname{Dom}\widetilde{\overline{\partial}}_b := \{ u \in L^2(\widetilde{m}) : \text{the distributional } \overline{\partial}_b \text{ of } u \text{ on } X \text{ is in } L^2_{(0,1)}(\widetilde{m}, \widehat{\theta}) \},\$$

and define

 $\widetilde{\overline{\partial}}_b u :=$ the distributional $\overline{\partial}_b$ of u on X

if $u \in \text{Dom}\widetilde{\overline{\partial}}_b$. Then

$$\widetilde{\overline{\partial}}_b$$
: Dom $\widetilde{\overline{\partial}}_b \subset L^2(\widetilde{m}) \to L^2_{(0,1)}(\widetilde{m}, \widehat{\theta})$

is a densely defined closed linear operator. Let

$$\tilde{\overline{\partial}}_b^*: \operatorname{Dom} \tilde{\overline{\partial}}_b^* \subset L^2_{(0,1)}(\tilde{m}, \hat{\theta}) \to L^2(\tilde{m})$$

be its adjoint. Let $\tilde{\Box}_b$ denote the Gaffney extension of the Kohn Laplacian given by

$$\operatorname{Dom} \tilde{\Box}_b = \{ s \in L^2(\tilde{m}) : s \in \operatorname{Dom} \tilde{\overline{\partial}}_b, \ \tilde{\overline{\partial}}_b s \in \operatorname{Dom} \tilde{\overline{\partial}}_b^* \}, \\ \tilde{\Box}_b s = \tilde{\overline{\partial}}_b^* \tilde{\overline{\partial}}_b s \quad \text{for } s \in \operatorname{Dom} \tilde{\Box}_b.$$

It is then a non-negative self-adjoint operator on $L^2(\tilde{m})$. The analysis of $\tilde{\Box}_b$ is not as well-understood, since this Kohn Laplacian (in particular, the operator $\tilde{\partial}_b^*$) is defined with respect to a non-smooth measure \tilde{m} . Nonetheless, it is this Kohn Laplacian that can be related to our operator of interest, namely $\Box_{b,1}$, in a simple manner. We will prove that since $\tilde{m} = |\psi|^2 m_1$,

$$u \in \text{Dom}\,\Box_{b,1} \text{ if and only if } \frac{u}{w} \in \text{Dom}\,\widetilde{\Box}_b,$$
(1.6)

$$\Box_{b,1} u = \psi \widetilde{\Box}_b(\psi^{-1} u), \quad \forall u \in \text{Dom} \,\Box_{b,1}.$$
(1.7)

Thus we can understand the solutions of $\Box_{b,1}$, once we understand the solutions of $\widetilde{\Box}_b$. In order to carry out the latter, we relate $\widetilde{\Box}_b$ to $\widehat{\Box}_b$: we will show that

$$\operatorname{Dom} \tilde{\Box}_b = \operatorname{Dom} \hat{\Box}_b, \tag{1.8}$$

and there exists some $g \in \mathcal{E}(\hat{\rho}^1, T^{0,1}\hat{X})$ (possibly non-smooth at p) such that

$$\widetilde{\Box}_b u = \hat{\Box}_b u + g \hat{\overline{\partial}}_b u, \quad \forall u \in \text{Dom}\, \widetilde{\Box}_b.$$
(1.9)

(Here $g\hat{\bar{\partial}}_b u$ is the pointwise pairing of the (0,1) vector g with the (0,1) form $\hat{\bar{\partial}}_b u$; see the discussion in Section 1.5 for the precise meaning of $\mathcal{E}(\hat{\rho}^1, T^{0,1}\hat{X})$.)

We can now outline the proofs of Theorems 1.3, 1.4 and 1.5.

First, from (1.6) and (1.7), it is clear that $\Box_{b,1}$ has closed range in $L^2(m_1)$, if and only if $\widetilde{\Box}_b$ has closed range in $L^2(\tilde{m})$. On the other hand, one can check that $\tilde{\bar{\partial}}_b$: Dom $\tilde{\bar{\partial}}_b \subset L^2(\tilde{m}) \to L^2_{(0,1)}(\tilde{m},\hat{\theta})$ is the identical as an operator to $\hat{\bar{\partial}}_b$: Dom $\hat{\bar{\partial}}_b \subset$ $L^2(\hat{m}) \to L^2_{(0,1)}(\hat{m},\hat{\theta})$. The latter is known to have closed range in $L^2(\hat{m})$ by the CR embeddability of \hat{X} ; see [17] (also [4]). Hence the same holds for the former, and it follows that $\tilde{\Box}_b$ has closed range in $L^2(\tilde{m})$. This proves Theorem 1.3.

Now from the above argument, we see that not only $\tilde{\Box}_b$ has closed range in $L^2(\tilde{m})$, but also $\hat{\Box}_b$ has closed range in $L^2(\hat{m})$. Thus there exist partial inverses

$$\tilde{N}: L^2(\tilde{m}) \to \text{Dom}(\tilde{\Box}_b) \subset L^2(\tilde{m}),$$

and

$$\hat{N}: L^2(\hat{m}) \to \text{Dom}(\hat{\Box}_b) \subset L^2(\hat{m})$$

to $\tilde{\Box}_b$ and $\hat{\Box}_b$ respectively, so that if

$$\tilde{\Pi}: L^2(\tilde{m}) \to L^2(\tilde{m})$$

is the orthogonal projection onto the kernel of $\tilde{\Box}_b$ in $L^2(\tilde{m})$, and

$$\hat{\Pi}: L^2(\hat{m}) \to L^2(\hat{m})$$

is the orthogonal projection onto the kernel of $\hat{\Box}_b$ in $L^2(\hat{m})$, then

$$\widetilde{\Box}_b \widetilde{N} + \widetilde{\Pi} = I \quad \text{on } L^2(\widetilde{m}),$$

and

$$\hat{\Box}_b \hat{N} + \hat{\Pi} = I \quad \text{on } L^2(\hat{m}).$$

From the relation

 $m_1 = |\psi|^{-2}\tilde{m}$

and (1.7), it is easy to see that

$$\Pi = \psi \widetilde{\Pi} \psi^{-1} \quad \text{and} \quad N = \psi \widetilde{N} \psi^{-1}, \tag{1.10}$$

at least when applied to functions in $L^2(m_1)$; thus to prove Theorem 1.4, it suffices to prove instead that $\widetilde{\Pi}$ and \widetilde{N} extend as continuous operators

$$\widetilde{\Pi}: \mathcal{E}(\hat{\rho}^{-4+\delta}) \to \mathcal{E}(\hat{\rho}^{-4+\delta})$$
(1.11)

$$\widetilde{N}: \mathcal{E}(\hat{\rho}^{-4+\delta}) \to \mathcal{E}(\hat{\rho}^{-2+\delta})$$
(1.12)

for every $0 < \delta < 2$. In order to do so, we relate $\widetilde{\Pi}$ to $\widehat{\Pi}$, and \widetilde{N} to \hat{N} , since $\widehat{\Pi}$ and \hat{N} are much better understood. We will show, on $L^2(\widetilde{m})$, that

$$\widetilde{\Pi}(I+\hat{R}) = \hat{\Pi} \tag{1.13}$$

$$\widetilde{N}(I+\hat{R}) = (I-\widetilde{\Pi})\hat{N}, \qquad (1.14)$$

where

$$\hat{R} := g \overline{\bar{\partial}}_b \hat{N} \colon L^2(\tilde{m}) \to L^2(\tilde{m})$$

is a continuous linear operator; in fact, these identities are almost immediate from (1.8) and (1.9). Furthermore, one can show that \hat{N} and $\hat{\Pi}$ extend as continuous operators

$$\hat{\Pi}: \mathcal{E}(\hat{\rho}^{-4+\delta}) \to \mathcal{E}(\hat{\rho}^{-4+\delta})$$
(1.15)

$$\hat{N}: \mathcal{E}(\hat{\rho}^{-4+\delta}) \to \mathcal{E}(\hat{\rho}^{-2+\delta})$$
(1.16)

for every $0 < \delta < 2$; thus if one can show that $(I + \hat{R})$ is invertible on $L^2(\tilde{m})$, and that the inverse extends to a continuous operator $\mathcal{E}(\hat{\rho}^{-4+\delta}) \to \mathcal{E}(\hat{\rho}^{-4+\delta})$, then from (1.13), at least one can conclude the assertion about $\widetilde{\Pi}$ in Theorem 1.4. It turns out that it is unclear whether or not the latter can be done; so we choose to proceed differently, by some bootstrap argument. It is this argument that we explain below.

First, so far $\widetilde{\Pi}$ and \widetilde{N} are defined only on $L^2(\widetilde{m})$. Since $\mathcal{E}(\hat{\rho}^{-4+\delta})$ is not a subset of $L^2(\widetilde{m})$ when $0 < \delta < 2$, we need to first extend $\widetilde{\Pi}$ and \widetilde{N} to $\mathcal{E}(\hat{\rho}^{-4+\delta})$, $0 < \delta < 2$. This is done by rewriting (1.13) and (1.14) as

$$\widetilde{\Pi} = \widehat{\Pi} - \widetilde{\Pi}\hat{R}, \quad \widetilde{N} = (I - \widetilde{\Pi})\hat{N} - \widetilde{N}\hat{R}.$$
(1.17)

Note that \hat{R} extends to a continuous operator

$$\hat{R}: \mathcal{E}(\hat{\rho}^{-4+\delta}) \to \mathcal{E}(\hat{\rho}^{-2+\delta}) \subset L^2(\widetilde{m})$$

for every $0 < \delta < 2$, since $\hat{R} = g\hat{\bar{\partial}}_b \hat{N}$, and \hat{N} satisfies the analogous property. Thus the second term on the right hand sides of the equations in (1.17) map $\mathcal{E}(\hat{\rho}^{-4+\delta})$ continuously to $L^2(\tilde{m})$. It follows that the domains of definition of Π and \tilde{N} can be extended to $\mathcal{E}(\hat{\rho}^{-4+\delta}), 0 < \delta < 2$.

To proceed further, let's write $\hat{\Pi}^{*,\widetilde{m}}$, $\hat{N}^{*,\widetilde{m}}$ and $\hat{R}^{*,\widetilde{m}}$ for the adjoints of $\hat{\Pi}$, \hat{N} and \hat{R} with respect to the inner product of $L^2(\widetilde{m})$. Since $\widetilde{\Pi}$ and \widetilde{N} are self-adjoint operators on $L^2(\widetilde{m})$, we have, by (1.13) and (1.14), that

$$(I + \hat{R}^{*,\widetilde{m}})\widetilde{\Pi} = \hat{\Pi}^{*,\widetilde{m}}$$
(1.18)

$$(I + \hat{R}^{*,\widetilde{m}})\tilde{N} = \hat{N}^{*,\widetilde{m}}(I - \widetilde{\Pi})$$
(1.19)

on $L^2(\widetilde{m})$. Now we need to understand some mapping properties of $\hat{\Pi}^{*,\widetilde{m}}$, $\hat{N}^{*,\widetilde{m}}$ and $\hat{R}^{*,\widetilde{m}}$; to do so, we note that

$$\begin{split} \hat{\Pi}^{*,\widetilde{m}} &= \frac{\hat{m}}{\widetilde{m}} \hat{\Pi} \frac{\widetilde{m}}{\hat{m}}, \\ \hat{N}^{*,\widetilde{m}} &= \frac{\hat{m}}{\widetilde{m}} \hat{N} \frac{\widetilde{m}}{\hat{m}}, \\ \hat{R}^{*,\widetilde{m}} &= \frac{\hat{m}}{\widetilde{m}} \hat{N} \hat{\bar{\partial}}_{b}^{*} (g^{*} \frac{\widetilde{m}}{\hat{m}}), \end{split}$$

where $\tilde{m}/\hat{m} := G_p^2 |\psi|^2$ is the density of \tilde{m} with respect to \hat{m} , and similarly $\hat{m}/\tilde{m} := G_p^{-2} |\psi|^{-2}$. Here g^* is the (0,1) form dual to g. Note that $\tilde{m}/\hat{m}, \hat{m}/\tilde{m} \in \mathcal{E}(\hat{\rho}^0)$. Hence one can show that

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$$\hat{\Pi}^{*,\widetilde{m}}: \mathcal{E}(\hat{\rho}^{-4+\delta}) \to \mathcal{E}(\hat{\rho}^{-4+\delta}) \tag{1.20}$$

$$\hat{N}^{*,\widetilde{m}}: \mathcal{E}(\hat{\rho}^{-4+\delta}) \to \mathcal{E}(\hat{\rho}^{-2+\delta})$$
(1.21)

$$\hat{R}^{*,\widetilde{m}}: \mathcal{E}(\hat{\rho}^{-4+\delta}) \to \mathcal{E}(\hat{\rho}^{-2+\delta})$$
(1.22)

for every $0 < \delta < 2$; these are easy consequences of the analogous properties of $\hat{\Pi}$, \hat{N} and \hat{R} . It then follows that (1.18) and (1.19) continue to hold on $\mathcal{E}(\hat{\rho}^{-4+\delta})$ for all $0 < \delta < 2$.

Now we return to (1.18) and (1.19). The problem facing us there is that we do not know whether $I + \hat{R}^{*,\widetilde{m}}$ is invertible on $\mathcal{E}(\hat{\rho}^{-4+\delta})$; if it is, then we can conclude, say from (1.18), that at least $\widetilde{\Pi}$ satisfies the conclusion of Theorem 1.4. In order to get around this problem, we have to proceed differently; the trick here is to introduce a suitable cut-off function, as follows.

Let χ be a smooth function on \hat{X} , such that χ is identically 1 in a neighborhood of p, and vanishes outside a small neighborhood of p. Then by (1.18) and (1.19), we have

$$(I + \hat{R}^{*,\widetilde{m}}\chi)\widetilde{\Pi} = \hat{\Pi}^{*,\widetilde{m}} - \hat{R}^{*,\widetilde{m}}(1-\chi)\widetilde{\Pi},$$
(1.23)

$$(I + \hat{R}^{*,\widetilde{m}}\chi)\tilde{N} = \hat{N}^{*,\widetilde{m}}(I - \widetilde{\Pi}) - \hat{R}^{*,\widetilde{m}}(1 - \chi)\tilde{N}$$
(1.24)

on $\mathcal{E}(\hat{\rho}^{-4+\delta})$ for all $0 < \delta < 2$. The upshot here is the following: if the support of χ is sufficiently small, then $(I + \hat{R}^{*,\widetilde{m}}\chi)$ is invertible on $L^2(\widetilde{m})$, and extends to a linear map

$$(I + \hat{R}^{*,\widetilde{m}}\chi)^{-1} \colon \mathcal{E}(\hat{\rho}^{-4+\delta}) \to \mathcal{E}(\hat{\rho}^{-4+\delta})$$
(1.25)

for every $0 < \delta < 4$. Roughly speaking, this is possible, because

$$I + \hat{R}^{*,\widetilde{m}}\chi = I + \frac{\hat{m}}{\widetilde{m}}\hat{N}\hat{\bar{\partial}}_{b}^{*}(\chi g^{*}\frac{\widetilde{m}}{\hat{m}}),$$

and because

$$\chi g^* \in \mathcal{E}(\hat{\rho}^1, \Lambda^{0,1}T^*\hat{X})$$

has compact support in a sufficiently small neighborhood of p. (In particular, χg^* is small.) Furthermore, for any $0 < \delta < 2$, one can show that

$$(1-\chi)\widetilde{\Pi}: \mathcal{E}(\hat{\rho}^{-4+\delta}) \to C_0^{\infty}(X)$$
(1.26)

$$(1-\chi)\widetilde{N}: \mathcal{E}(\hat{\rho}^{-4+\delta}) \to C_0^{\infty}(X)$$
(1.27)

where $C_0^{\infty}(X)$ is the space of all smooth functions on X that has compact support in X. These can be used to control the last term on the right hand side of (1.23) and (1.24). By (1.20) and (1.22), one then concludes that the right hand side of (1.23) maps $\mathcal{E}(\hat{\rho}^{-4+\delta})$ into itself for $0 < \delta < 2$; thus (1.25) shows that (1.11) holds as desired. This in turn controls the first term of (1.24); by (1.21), (1.22) and (1.27), one concludes that the right hand side of (1.24) maps $\mathcal{E}(\hat{\rho}^{-4+\delta})$ into $\mathcal{E}(\hat{\rho}^{-2+\delta})$ for $0 < \delta < 2$. Finally, another application of (1.25) shows that \tilde{N} satisfies (1.12) as desired. Thus Theorem 1.4 is established, modulo (1.25), (1.26) and (1.27).

It may help to reiterate here the reason for the introduction of the cut-off χ : that was introduced so that one can invert $I + \hat{R}^{*,\widetilde{m}}\chi$. In fact, since the coefficient of $g^* \in \mathcal{E}(\hat{\rho}^1, \Lambda^{0,1}T^*\hat{X}), \chi g^*$ has sufficiently small L^{∞} norm, if the support of χ is chosen sufficiently small. As a result, one could make $\|\hat{R}^{*,\widetilde{m}}\chi\|_{L^2(\widetilde{m})\to L^2(\widetilde{m})} \leq 1/2$, by controlling the support of χ . This allows one to invert $I + \hat{R}^{*,\widetilde{m}}\chi$ on $L^2(\widetilde{m})$ via a Neumann series. We note in passing that it is only because we have taken the adjoint of \hat{R} that the product χg^* appears in the expression for $I + \hat{R}^{*,\widetilde{m}}\chi$; that is essentially why we want to take the adjoints of (1.13) and (1.14). One can then proceed to extend $(I + \hat{R}^{*,\widetilde{m}}\chi)^{-1}$ so that it satisfies (1.25); the precise detail is rather involved, and we leave this until Section 3.

It then remains to prove (1.26) and (1.27). To do so, we need yet to introduce yet another Kohn Laplacian, namely $\hat{\Box}_{b,\varepsilon}$. Let $\eta \in C_0^{\infty}(\mathbb{R}^3)$ be a non-negative function such that $\eta(z,t) = 1$ if $\hat{\rho}(z,t) \leq 1/2$, and $\eta(z,t) = 0$ if $\hat{\rho}(z,t) \geq 1$. For $0 < \varepsilon < 1$, let $\eta_{\varepsilon}(z,t) = \eta(\varepsilon^{-1}z, \varepsilon^{-2}t)$, and let

$$\hat{m}_{\varepsilon} := \eta_{\varepsilon} \hat{m} + (1 - \eta_{\varepsilon}) \widetilde{m}.$$

This volume form has a smooth density against \hat{m} , and the volume of \hat{X} with respect to this volume form is finite. So if we extend the Cauchy–Riemann operator such that

$$\operatorname{Dom}\bar{\bar{\partial}}_{b,\varepsilon} := \{ u \in L^2(\hat{m}_{\varepsilon}) : \text{the distributional } \overline{\partial}_b \text{ of } u \text{ is in } L^2_{(0,1)}(\hat{m}_{\varepsilon},\hat{\theta}) \},\$$

and define

$$\hat{\overline{\partial}}_{b,\varepsilon} u :=$$
the distributional $\overline{\partial}_b$ of u

if $u \in \text{Dom}\overline{\overline{\partial}}_{b,\varepsilon}$, then

$$\hat{\bar{\partial}}_{b,\varepsilon}$$
: Dom $\hat{\bar{\partial}}_{b,\varepsilon} \subset L^2(\hat{m}_{\varepsilon}) \to L^2_{(0,1)}(\hat{m}_{\varepsilon},\hat{\theta})$

is a densely defined closed linear operator. Let

$$\hat{\bar{\partial}}_{b,\varepsilon}^*: \operatorname{Dom} \hat{\bar{\partial}}_{b,\varepsilon}^* \subset L^2_{(0,1)}(\hat{m}_{\varepsilon}, \hat{\theta}) \to L^2(\hat{m}_{\varepsilon})$$

be its adjoint. Let $\hat{\Box}_{b,\varepsilon}$ be the Gaffney extension of the Kohn Laplacian $\hat{\overline{\partial}}_{b,\varepsilon}^* \hat{\overline{\partial}}_{b,\varepsilon}$. Then $\hat{\Box}_{b,\varepsilon}$ is almost as well-behaved as $\hat{\Box}_b$. In particular, $\hat{\Box}_{b,\varepsilon}$: Dom $\hat{\Box}_{b,\varepsilon} \subset L^2(\hat{m}_{\varepsilon}) \to L^2(\hat{m}_{\varepsilon})$

$$\hat{\Box}_{b,\varepsilon}\hat{N}_{\varepsilon} + \hat{\Pi}_{\varepsilon} = I.$$

Furthermore,

$$\operatorname{Dom} \tilde{\Box}_b = \operatorname{Dom} \hat{\Box}_{b,\varepsilon},$$

and there will exist some $g_{\varepsilon} \in \mathcal{E}(\hat{\rho}^1, T^{0,1}\hat{X})$ (possibly non-smooth near p) such that

$$\widetilde{\Box}_{b}u = \hat{\Box}_{b,\varepsilon}u + g_{\varepsilon}\hat{\overline{\partial}}_{b}u, \quad \forall u \in \text{Dom}\,\tilde{\Box}_{b}.$$
(1.28)

The upshot here is that g_{ε} will be compactly supported in the support of η_{ε} , whereas our previous g may not be compactly supported near p. One can repeat the proof of (1.18) and (1.19), and show that

$$\begin{split} (I + \hat{R}_{\varepsilon}^{*,\widetilde{m}}) \widetilde{\Pi} &= \hat{\Pi}_{\varepsilon}^{*,\widetilde{m}} \\ (I + \hat{R}_{\varepsilon}^{*,\widetilde{m}}) \widetilde{N} &= \hat{N}_{\varepsilon}^{*,\widetilde{m}} (I - \widetilde{\Pi}) \end{split}$$

on $L^2(\widetilde{m})$, where

$$\hat{R}_{\varepsilon} := g_{\varepsilon} \overline{\overline{\partial}}_{b,\varepsilon} \hat{N}_{\varepsilon}$$

It follows that

$$(1-\chi)\widetilde{\Pi} = (1-\chi)\widehat{\Pi}_{\varepsilon}^{*,\widetilde{m}} - (1-\chi)\widehat{R}_{\varepsilon}^{*,\widetilde{m}}\widetilde{\Pi}, \qquad (1.29)$$

$$(1-\chi)\widetilde{N} = (1-\chi)\hat{N}_{\varepsilon}^{*,\widetilde{m}}(I-\widetilde{\Pi}) - (1-\chi)\hat{R}_{\varepsilon}^{*,\widetilde{m}}\widetilde{N}.$$
(1.30)

Now

$$(1-\chi)\hat{R}^{*,\widetilde{m}}_{\varepsilon} = \frac{\hat{m}_{\varepsilon}}{\widetilde{m}}(1-\chi)\hat{N}_{\varepsilon}\hat{\overline{\partial}}^{*}_{b,\varepsilon}g^{*}_{\varepsilon}\frac{\widetilde{m}}{\hat{m}_{\varepsilon}},$$

and if ε is chosen sufficiently small (so that the support of g_{ε} is disjoint from that of $1-\chi$), then $(1-\chi)\hat{N}_{\varepsilon}\hat{\overline{\partial}}_{b,\varepsilon}^{*}g_{\varepsilon}^{*}$ is an infinitely smoothing pseudodifferential operator, by pseudolocality of \hat{N}_{ε} . Hence the last term of (1.29), and also the last term of (1.30), map $\mathcal{E}(\hat{\rho}^{-4+\delta})$ into $C_{0}^{\infty}(X)$. Since for every $0 < \varepsilon < 1$ and every $0 < \delta < 2$,

$$\hat{\Pi}^{*,m}_{\varepsilon}: \mathcal{E}(\hat{\rho}^{-4+\delta}) \to \mathcal{E}(\hat{\rho}^{-4+\delta}),$$

and

$$\hat{N}_{\varepsilon}^{*,m}: \mathcal{E}(\hat{\rho}^{-4+\delta}) \to \mathcal{E}(\hat{\rho}^{-2+\delta}),$$

it follows that $(1 - \chi)\widetilde{\Pi}$ and $(1 - \chi)\widetilde{N}$ satisfy (1.26) and (1.27), and we are done with the proof of Theorem 1.4.

Finally, to prove Theorem 1.5, the key is the following fact, which we prove in Lemma 9.1: if α is a (0,1) form with coefficients in $\mathcal{E}(\hat{\rho}^0)$, then $\bar{\partial}_{b,1}^* \alpha \in \mathcal{E}(\hat{\rho}^{-1})$ satisfies $\Pi \bar{\partial}_{b,1}^* \alpha = 0$. To compute $\Pi f = \Pi \Box_{b,1} \tilde{\beta} = \Pi \bar{\partial}_{b,1}^* (\bar{\partial}_{b,1} \tilde{\beta})$, we will then decompose $\bar{\partial}_{b,1} \tilde{\beta}$ into a sum

$$\overline{\partial}_{b,1}\widetilde{\beta} = \alpha_0 + E,$$

where the main term α_0 has coefficients in $\mathcal{E}(\hat{\rho}^{-2})$, and the error E has coefficients in $\mathcal{E}(\hat{\rho}^0)$. Then by Lemma 9.1, $\Pi \overline{\partial}_{b,1}^* E = 0$. Furthermore, we will construct by hand an explicit family of (0, 1) forms α_{ε} , with coefficients in $\mathcal{E}(\hat{\rho}^0)$, such that

$$\overline{\partial}_{b,1}^* \alpha_{\varepsilon} \to \overline{\partial}_{b,1}^* \alpha_0 \quad \text{in } \mathcal{E}(\hat{\rho}^{-1}) \text{ as } \varepsilon \to 0.$$

Thus by continuity of Π on $\mathcal{E}(\hat{\rho}^{-1})$, we have $\Pi \overline{\partial}_{b,1}^* \alpha_0 = \lim_{\varepsilon \to 0} \Pi \overline{\partial}_{b,1}^* \alpha_{\varepsilon} = 0$ as well, the last equality following from Lemma 9.1. Together we get $\Pi f = 0$, as desired.

1.5. Definitions and notations

We shall now recall some basic definitions, and introduce some basic notations.

A 3-dimensional smooth manifold X is said to be a CR manifold, if there exists a 1-dimensional subbundle L of the complexified tangent bundle $\mathbb{C}TX$ such that $L \cap \overline{L} = \{0\}$; such subbundle L is then denoted as $T^{1,0}X$, and \overline{L} denoted as $T^{0,1}X$. The dual bundles to $T^{1,0}X$ and $T^{0,1}X$ will be denoted by $\Lambda^{1,0}T^*X$ and $\Lambda^{0,1}T^*X$ respectively. A typical example is a 3-dimensional smooth submanifold X of \mathbb{C}^N ; there one has a natural CR structure induced from \mathbb{C}^N , given by the bundle of all (1,0) vectors in \mathbb{C}^N that are tangent to X.

A 3-dimensional CR manifold X is said to be strongly pseudoconvex, if at every point on X there exists a local section Z of $T^{1,0}X$ such that $[Z,\overline{Z}]$ is transverse to $T^{1,0}X \oplus T^{0,1}X$. It is said to be CR embeddable in \mathbb{C}^N , if there exists a smooth embedding $\Phi: X \to \Phi(X) \subset \mathbb{C}^N$, such that $d\Phi(T^{1,0}X)$ agrees with the natural CR structure of $\Phi(X)$ induced from \mathbb{C}^N .

We shall write $C^{\infty}(X)$ for the space of smooth functions on X, and $\Omega_0^{0,1}(X)$ for the space of smooth sections of $\Lambda^{0,1}T^*X$. We shall also write $C_0^{\infty}(X)$ and $\Omega_0^{0,1}(X)$ for the subspaces of $C^{\infty}(X)$ and $\Omega^{0,1}(X)$ which consist of elements that have compact support in X.

Suppose X is a 3-dimensional CR manifold. If there exists a real contact form θ (i.e. a global real 1-form θ with $\theta \wedge d\theta \neq 0$ everywhere) such that $T^{1,0}X \oplus T^{0,1}X$ is given by the kernel of θ , then $(X, T^{1,0}X, \theta)$ is called a pseudohermitian 3-manifold. In that case, X is strongly pseudoconvex, and one can define a Hermitian inner product on $T^{1,0}X$, by

$$\langle Z_1 | Z_2 \rangle_{\theta} := \frac{1}{2} d\theta (Z_1, i \overline{Z}_2).$$

This allows one to define various geometric quantities on X, like the connection form ω_1^1 , and the Tanaka–Webster scalar curvature R. One can also define the sublaplacian Δ_b , the conformal sublaplacian $L_b := -4\Delta_b + R$, and the CR Paneitz operator P. We refer the reader to say [5] for the precise definitions of such.

We note that the above Hermitian inner product on $T^{1,0}X$ induces naturally a Hermitian inner product on $\Lambda^{0,1}T^*X$, which we still denote by $\langle \cdot | \cdot \rangle_{\theta}$. For $\alpha \in \Omega^{0,1}X$, we write $|\alpha|_{\theta}^2 := \langle \alpha | \alpha \rangle_{\theta}$.

In what follows, we will need the Reeb vector field T on a contact manifold (X, θ) , which is the unique vector field such that

$$\theta(T) \equiv 1, \quad d\theta(T, \cdot) \equiv 0.$$

If $\hat{\rho}$ is a non-negative smooth function defined near a point p, then we write $\varepsilon(\hat{\rho}^k)$ for the set of smooth functions $f \in C^{\infty}(\hat{X})$ such that $|f| \leq C\hat{\rho}^k$ near p for some C > 0.

Suppose now $(\hat{X}, T^{1,0}\hat{X}, \hat{\theta}_0)$ is a pseudohermitian 3-manifold, and we fix $p \in \hat{X}$. Then as is known, there exists another contact form $\hat{\theta}$ on \hat{X} , which is a multiple of $\hat{\theta}_0$ by a positive smooth function, so that near p, there exist CR normal coordinates (z, t). In other words, the contact form $\hat{\theta}$ and the coordinates (z, t) are chosen, so that

- (i) the point p corresponds to (z, t) = (0, 0);
- (ii) one can find a local section \hat{Z}_1 of $T^{1,0}\hat{X}$ near p, with $\langle \hat{Z}_1 | \hat{Z}_1 \rangle_{\hat{\theta}} = 1$, such that \hat{Z}_1 admits the following expansion near p:

$$\hat{Z}_1 = \frac{\partial}{\partial z} - i\overline{z}\frac{\partial}{\partial t} + \varepsilon(\hat{\rho}^4)\frac{\partial}{\partial z} + \varepsilon(\hat{\rho}^4)\frac{\partial}{\partial \overline{z}} + \varepsilon(\hat{\rho}^5)\frac{\partial}{\partial t}; \qquad (1.31)$$

(iii) the Reeb vector field \hat{T} with respect to $\hat{\theta}$ admits an expansion

$$\hat{T} = \frac{\partial}{\partial t} + \varepsilon(\hat{\rho}^3)\frac{\partial}{\partial z} + \varepsilon(\hat{\rho}^3)\frac{\partial}{\partial \overline{z}} + \varepsilon(\hat{\rho}^4)\frac{\partial}{\partial t}.$$
(1.32)

Here

$$\hat{\rho}(z,t) = (|z|^4 + t^2)^{\frac{1}{4}}$$

for (z,t) in a neighborhood of (0,0). For later convenience, from now on we will fix a positive smooth extension of $\hat{\rho}$ to the whole manifold X. We will also write $\hat{Z}_{\bar{1}} := \overline{\hat{Z}_1}$. Note that in CR normal coordinates we have

$$\hat{\theta} = dt - i(zd\bar{z} - \bar{z}dz) + \varepsilon(\hat{\rho}^5)dz + \varepsilon(\hat{\rho}^5)d\bar{z} + \varepsilon(\hat{\rho}^4)dt.$$
(1.33)

Next, for $m \in \mathbb{R}$, we will introduce a Fréchet space $\mathcal{E}(\hat{\rho}^m)$, with which our results are formulated. We pause and introduce some notations first. Let $k \in \mathbb{N}$. We denote by $\hat{\nabla}_b^k$ any differential operator of the form $L_1 \dots L_k$, where $L_j \in C^{\infty}(\hat{X}, T^{1,0}\hat{X} \oplus T^{0,1}\hat{X})$, $\langle L_j | L_j \rangle_{\hat{\theta}} \leq 1, j = 1, \dots, k$. Let $\mathcal{O}(\hat{\rho}^m) = \mathcal{O}^{(0)}(\hat{\rho}^m), m \in \mathbb{R}$, denote the set of all $f \in C^{\infty}(X)$ such that $|f| \leq C\hat{\rho}^m$ near p, for some C > 0. Let $\mathcal{O}^{(1)}(\hat{\rho}^m)$ denote the set of all functions $f \in \mathcal{O}(\hat{\rho}^m)$ such that $\hat{\nabla}_b f \in \mathcal{O}(\hat{\rho}^{m-1})$. Similarly, for $k \in \mathbb{N}, k \geq 2$, let $\mathcal{O}^{(k)}(\hat{\rho}^m)$ denote the set of all functions $f \in \mathcal{O}(\hat{\rho}^m)$ such that $\hat{\nabla}_b f \in \mathcal{O}(\hat{\rho}^m)$ such that $\hat{\nabla}_b f \in \mathcal{O}^{(k-1)}(\hat{\rho}^{m-1})$. Put

$$\mathcal{E}(\hat{\rho}^m) = \bigcap_{k \in \mathbb{N} \bigcup \{0\}} \mathcal{O}^{(k)}(\hat{\rho}^m).$$
(1.34)

Let $\Omega \subset \hat{X}$ be an open set. For $f \in C^{\infty}(\Omega)$, define

$$\|f\|_{L^{\infty}(\Omega)} := \sup_{x \in \Omega} |f(x)|.$$

 $\mathcal{E}(\hat{\rho}^m)$ is a Fréchet space with the semi-norms:

$$u \to \left\| \hat{\nabla}_b^k(\hat{\rho}^{-m+k}u) \right\|_{L^{\infty}(X)}, \quad u \in \mathcal{E}(\hat{\rho}^m), \tag{1.35}$$

for $k \in \mathbb{N}_0$. These semi-norms then define the topology of $\mathcal{E}(\hat{\rho}^m)$.

There is a version of this space for smooth vector bundles over X. Let E be a smooth vector bundle over \hat{X} of rank r. Let f_1, \ldots, f_r be any local frame in some small neighborhood U of p. For $m \in \mathbb{R}$, let $\mathcal{E}(\hat{\rho}^m, E)$ be the set of all $u \in C^{\infty}(X, E)$ such that $u = u_1 f_1 + \cdots + u_r f_r$ on U, and $\chi u_j \in \mathcal{E}(\hat{\rho}^m)$ for every $\chi \in C_0^{\infty}(U)$ and every $j = 1, \ldots, r$.

Note that $\mathcal{E}(\hat{\rho}^m) \subset \mathcal{E}(\hat{\rho}^{m'})$ if m' < m. We also notice that for every $m \in \mathbb{R}$, $C_0^{\infty}(X)$ is dense in $\mathcal{E}(\hat{\rho}^m)$ for the topology of $\mathcal{E}(\hat{\rho}^{m'})$, for every m' < m. Similarly for $\mathcal{E}(\hat{\rho}^m, E)$ for any smooth vector bundle E.

Distributions on \hat{X} will be denoted by $\mathcal{D}'(\hat{X})$.

Finally, suppose $T: \text{Dom}(T) \subset H_1 \to H_2$ is a densely defined closed linear operator between two Hilbert spaces H_1 and H_2 , and suppose the range of T is closed in H_2 . Then the partial inverse of T is the unique linear operator $S: H_2 \to H_1$, such that if $\Pi_1: H_1 \to H_1$ and $\Pi_2: H_2 \to H_2$ are the orthogonal projections onto the kernels of T and T^* respectively, then

$$TS + \Pi_2 = I$$
, $S\Pi_2 = 0$, and $\Pi_1 S = 0$.

(Here $T^*: \text{Dom}(T^*) \subset H_2 \to H_1$ is the adjoint of T, which is also densely defined and closed; and I is the identity operator.) It follows that $S: H_2 \to H_1$ is bounded, and

 $ST + \Pi_1 = I$ on the domain of T.

As we saw in Section 1.4, the operators

$$\begin{split} \hat{\bar{\partial}}_b: \mathrm{Dom} \hat{\bar{\partial}}_b \subset L^2(\hat{m}) \to L^2_{(0,1)}(\hat{m}), \\ \hat{\Box}_b: \mathrm{Dom} \hat{\Box}_b \subset L^2(\hat{m}) \to L^2(\hat{m}), \\ \tilde{\Box}_b: \mathrm{Dom} \tilde{\Box}_b \subset L^2(\tilde{m}) \to L^2(\tilde{m}), \end{split}$$

and

$$\Box_{b,1}: \text{Dom}\Box_{b,1} \subset L^2(m_1) \to L^2(m_1)$$

all have closed ranges. (See Section 7 for more details.) Their partial inverses will be denoted by \hat{K} , \hat{N} , \tilde{N} and N respectively.

A piece of convention here: recall that \hat{K} is an operator that takes (0,1) forms to functions. By identifying the space of (0,1) forms locally with functions, we would sometimes like to think of \hat{K} as a map from functions to functions. To do so rigorously, we proceed as follows. At every point $x \in \hat{X}$, there exists a non-isotropic ball $B(x, r_x)$ such that $T^{0,1}\hat{X}$ has a non-zero section on $B(x, r_x)$. The set $\{B(x, r_x/2): x \in \hat{X}\}$ covers \hat{X} ; one can thus take a finite subcover that covers \hat{X} . Denote this finite subcover by $\{B_1, \ldots, B_N\}$, and the dual of a non-zero local section of $T^{0,1}\hat{X}$ on $2B_i$ by $\hat{\omega}_i$; here $2B_i$ is the non-isotropic ball that has the same center as B_i , but twice the radius. We further normalize $\hat{\omega}_i$ so that $\langle \hat{\omega}_i | \hat{\omega}_i \rangle_{\hat{\theta}} = 1$ on $2B_i$. One then has the following property: there is some $r_0 > 0$ such that if B is a non-isotropic ball of radius $< r_0$ on \hat{X} that intersects some of the B_i above, then $\hat{\omega}_i$ is defined and of norm 1 on B. By taking a partition of unity $\sum \eta_i = 1$ subordinate to the open cover $\{B_1, \ldots, B_N\}$, we can define maps \hat{K}_i , $1 \leq i \leq N$, that map from functions to functions, by the following formula:

$$\hat{K}_i \varphi := \hat{K}(\eta_i \varphi \hat{\omega}_i). \tag{1.36}$$

Then since $\hat{K} = \sum_{i} \hat{K} \eta_{i}$, and $\eta_{i} \phi = \eta_{i} \langle \phi | \hat{\omega}_{i} \rangle_{\hat{\theta}} \hat{\omega}_{i}$ for all (0, 1) forms ϕ , we have

$$\hat{K}\phi = \sum_{i} \hat{K}_{i} [\langle \phi | \hat{\omega}_{i} \rangle_{\hat{\theta}}].$$

It will be slightly more convenient to consider properties of \hat{K}_i instead of \hat{K} at a number of places below. The result will always be independent (up to constants) of the choices of the cut-offs η_i , and of the choice of frames $\hat{\omega}_i$.

The plan of the paper is as follows. In Section 2, we gather together some properties of the Szegö projection $\hat{\Pi}$, as well as the partial inverse \hat{N} of the smooth Kohn Laplacian $\hat{\Box}_b$. In Section 3, we develop tools to establish the key mapping property (1.25), that involves the weighted space $\mathcal{E}(\hat{\rho}^{-4+\delta})$. In Section 4, we construct the CR function ψ that is crucial for us. Sections 5 and 6 clarify the relations between the various Kohn Laplacians. Sections 7 to 9 contain the proofs of Theorems 1.3 to 1.5, which implies Theorem 1.1 and Corollary 1.2 as we have explained above. Finally, in Appendix A, we establish some properties of the Green's function of the conformal Laplacian L_b , which allows us to apply our results towards the study of the CR positive mass theorem as was laid out in [5]. In Appendix B, we prove a subelliptic estimate for $\hat{\Box}_b$, which should be known to the experts, but which has not appeared explicitly in literature.

2. Some properties of the smooth Kohn Laplacian $\hat{\Box}_b$

We collect in this section some results from subelliptic analysis and several complex variables. The key is to introduce a class of non-isotropic smoothing operators on our pseudohermitian manifold \hat{X} , and show that the Szegö projection $\hat{\Pi}$, as well as the partial inverse \hat{N} of $\hat{\Box}_b$, are examples of such; we will deduce, as a result, mapping properties of $\hat{\Pi}$ and \hat{N} with respect to the weighted spaces $\mathcal{E}(\hat{\rho}^{\delta})$. Many of these are known; we refer the reader to Nagel, Stein and Wainger [23], Kohn [16], Christ [6,7], Nagel, Rosay, Stein and Wainger [25], Koenig [15] and the references therein for further details.

First, we recall the Carnot–Caratheodory metric on our CR manifold \hat{X} . For $\delta > 0$, let $C(\delta)$ be the class of all absolutely continuous mappings $\varphi: [0, 1] \to \hat{X}$ such that for a.e. t,

$$\varphi'(t) = a_1(t)X_1(\varphi(t)) + a_2(t)X_2(\varphi(t)), \quad |a_j(t)| < \delta, \quad j = 1, 2.$$

Here X_1 and X_2 are the real and imaginary parts of \hat{Z}_1 respectively. The Carnot–Caratheodory metric on \hat{X} is then defined by

 $\vartheta(x,y) = \inf\{\delta > 0: \text{ there exists } \varphi \in C(\delta) \text{ such that } \varphi(0) = x, \varphi(1) = y\}$

for $x, y \in \hat{X}$. From Theorem 4 of [23], coupled with the representations (1.31) and (1.32) of \hat{Z}_1 and \hat{T} in CR normal coordinates, it is easy to show that for points x sufficiently close to p, we have

$$\vartheta(x,p) \simeq \hat{\rho}(x).$$

(See also Theorem 3.5 and Remark 3.3 of Jean [14].) We write B(x, r) for the nonisotropic ball $\{y \in \hat{X} : \vartheta(x, y) < r\}$ of radius r centered at x.

Next, we proceed to define on \hat{X} a class of (non-isotropic) smoothing operators of order j. For our purposes, it suffices to restrict our attention to the case when $0 \le j < 4$.

Recall that a function ϕ on \hat{X} is said to be a *normalized bump function* on a ball B(x,r), if it is smooth with compact support on B(x,r), and satisfies

$$\left\|\hat{\nabla}_{b}^{k}\phi\right\|_{L^{\infty}(B(x,r))} \le C_{k}r^{-k} \tag{2.1}$$

for all $k \ge 0$; here $C_k > 0$ are absolute constants independent of r.

Usually we only require the above derivative estimate to be satisfied for all $0 \le k \le N$ for some large integer N. In that case, we say that ϕ is a normalized bump function of order N in B(x, r).

Suppose now T is a continuous linear operator $T: C^{\infty}(\hat{X}) \to C^{\infty}(\hat{X})$, and its adjoint T^* (with respect to the inner product of $L^2(\hat{m})$) is also a continuous map $T^*: C^{\infty}(\hat{X}) \to C^{\infty}(\hat{X})$. We say that T is a smoothing operator of order $j, 0 \leq j < 4$, if

(a) there exists a kernel T(x, y), defined and smooth away from the diagonal in $\hat{X} \times \hat{X}$, such that

$$Tf(x) = \int_{\hat{X}} T(x,y)f(y)\hat{m}(y)$$
(2.2)

for any $f \in C^{\infty}(\hat{X})$, and every x not in the support of f;

(b) the kernel T(x, y) satisfies the following differential inequalities when $x \neq y$:

$$|(\hat{\nabla}_b)_x^{\alpha_1}(\hat{\nabla}_b)_y^{\alpha_2}T(x,y)| \lesssim_{\alpha} \vartheta(x,y)^{-4+j-|\alpha|}, \quad |\alpha| = |\alpha_1| + |\alpha_2|;$$

(c) the operators T and T^* satisfy the following cancellation conditions: if ϕ is a normalized bump function in some ball B(x, r), then

$$\|\hat{\nabla}^{\alpha}_{b}T\phi\|_{L^{\infty}(B(x,r))} \lesssim_{\alpha} r^{j-|\alpha|},$$

and

$$\|\hat{\nabla}^{\alpha}_{b}T^{*}\phi\|_{L^{\infty}(B(x,r))} \lesssim_{\alpha} r^{j-|\alpha|}.$$

It is then clear that T is smoothing of order j, if and only if T^* is smoothing of order j. We also have the following proposition:

Proposition 2.1. If T is a smoothing operator of order 0, then T is bounded on $L^p(\hat{m})$ for 1 .

Proof. The boundedness of T on $L^2(\hat{m})$ follows from a version of T(1) theorem. In fact, suppose f is a normalized bump function on a ball $B(x_0, r)$. If T is a smoothing operator of order 0, then by the cancellation condition on T,

$$||Tf||_{L^2(B(x_0,2r))} \lesssim r^2,$$

and by the kernel representation of T, when $x \notin B(x_0, 2r)$,

$$|Tf(x)| \lesssim \int_{y \in B(x_0,r)} \frac{|f(y)|}{\vartheta(x,y)^4} \hat{m}(y) \lesssim \vartheta(x,x_0)^{-4} r^4.$$

Hence

$$\int_{\substack{x \notin B(x_0,2r)}} |Tf(x)|^2 \hat{m}(x) \lesssim \int_{\substack{x \notin B(x_0,2r)}} \vartheta(x,x_0)^{-8} r^8 \hat{m}(x) \lesssim r^4.$$

Altogether,

$$\|Tf\|_{L^2(\hat{X})} \lesssim r^2;$$

similarly for $||T^*f||_{L^2(\hat{X})}$. Hence both T and T^* are restrictedly bounded, and by the T(1) theorem (see e.g. Chapter 7 of [27]), T is bounded on $L^2(\hat{m})$. By the Calderon–Zygmund theory of singular integrals, it then follows that such operators are bounded on $L^p(\hat{m})$, $1 . (See e.g. Chapter 1 of [27].) <math>\Box$

Next we have the following theorem:

Theorem 2.2. The Szegö projection $\hat{\Pi}$, and the partial inverse \hat{N} of $\hat{\Box}_b$, are smoothing operators of orders 0 and 2 respectively. Furthermore, \hat{K} is smoothing of order 1, in the sense that the local representations \hat{K}_i defined by (1.36) are smoothing of order 1.

We defer its proof until the end of this section.

We will need two further key facts about this class of smoothing operators:

Theorem 2.3. If T_1 and T_2 are smoothing operators of orders j_1 and j_2 respectively, with $j_1, j_2 \ge 0$ and $j_1 + j_2 < 4$, then $T_1 \circ T_2$ is a smoothing operator of order $j_1 + j_2$.

Theorem 2.4. If T is a smoothing operator of order j, $0 \le j < 4$, then T extends to a continuous linear map

$$T: \mathcal{E}(\hat{\rho}^{-\gamma}) \to \mathcal{E}(\hat{\rho}^{-\gamma+j}),$$

as long as $j < \gamma < 4$.

In particular, in proving Theorem 2.2, it suffices to prove the statements for $\hat{\Pi}$ and \hat{K} , since the statement for \hat{N} follows from Theorem 2.3 and the well-known fact that

$$\hat{N} = \hat{K}\hat{K}^*$$

Also, combining Theorems 2.2 and 2.4, $\hat{\Pi}$ and \hat{N} extend to continuous linear maps

$$\begin{split} \hat{\Pi} &: \mathcal{E}(\hat{\rho}^{-4+\delta}) \to \mathcal{E}(\hat{\rho}^{-4+\delta}), \quad 0 < \delta < 4, \\ \hat{N} &: \mathcal{E}(\hat{\rho}^{-4+\delta}) \to \mathcal{E}(\hat{\rho}^{-2+\delta}), \quad 0 < \delta < 2, \end{split}$$

and (1.15), (1.16) follow.

Proof of Theorem 2.3. We will only need the case when T_1 is smoothing of order 1, and T_2 is smoothing of order 0 or 1. Thus we will focus on these cases.

Suppose first both T_1 and T_2 are smoothing of order 1. Then $T := T_1 \circ T_2$ is continuous on $C^{\infty}(\hat{X})$, and so is T^* . Furthermore, when $f \in C^{\infty}(\hat{X})$, we have

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$$T_j f(x) = \int_{\hat{X}} T_j(x, y) f(y) \hat{m}(y), \quad j = 1, 2$$

for all x not in the support of f. Using the cancellation conditions, one can show that this integral representation actually holds for all $x \in \hat{X}$ (not just for all x outside the support of f; cf. Chapter 1.7 of Stein [27]). This is typical of operators that are smoothing of positive orders.

As a result, by Fubini's theorem, when $f \in C^{\infty}(\hat{X})$, (2.2) holds for all $x \in \hat{X}$, where

$$T(x,y) = \int_{\hat{X}} T_1(x,z) T_2(z,y) \hat{m}(z).$$
(2.3)

Fix now $x, y \in \hat{X}$, and let $r = \vartheta(x, y)/4$. We pick normalized bump functions χ_1, χ_2 in B(x, r) and B(y, r) respectively, such that $\chi_1 \equiv 1$ on B(x, r/2), and $\chi_2 \equiv 1$ on B(y, r/2). Then inserting $1 = \chi_1(z) + \chi_2(z) + (1 - \chi_1 - \chi_2)(z)$ into the integral defining T(x, y), we have

$$T(x,y) = T^{(1)}(x,y) + T^{(2)}(x,y) + T^{(3)}(x,y),$$

and we estimate these one by one.

First,

$$T^{(1)}(x,y) = \int_{\hat{X}} T_1(x,z)\chi_1(z)T_2(z,y)\hat{m}(z).$$

We can differentiate under the integral, and obtain

$$(\hat{\nabla}_b)_y^{\alpha_2} T^{(1)}(x,y) = (T_1 f_y^{(\alpha_2)})(x)$$

where

$$f_y^{(\alpha_2)}(z) := \chi_1(z) (\hat{\nabla}_b)_y^{\alpha_2} T_2(z, y)$$

is $r^{-3-|\alpha_2|}$ times a normalized bump function in B(x,r). Thus by the cancellation condition for T_1 , we obtain

$$|(\hat{\nabla}_b)_x^{\alpha_1}(\hat{\nabla}_b)_y^{\alpha_2}T^{(1)}(x,y)| \lesssim r^{-3-|\alpha_2|}r^{1-|\alpha_1|} = r^{-2-|\alpha|}.$$

This proves the desired differential inequalities for $T^{(1)}(x, y)$. A similar argument, using the cancellation conditions for \overline{T}_2^* instead, shows that

$$|(\hat{\nabla}_b)_x^{\alpha_1}(\hat{\nabla}_b)_y^{\alpha_2}T^{(2)}(x,y)| \lesssim r^{-2-|\alpha|}.$$

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Finally, the integral defining $T^{(3)}(x, y)$ is supported for z outside the balls B(x, r/2)and B(y, r/2). As a result, $\vartheta(x, z) \simeq \vartheta(y, z)$ for z in the support of the integral defining $T^{(3)}(x, y)$. One can now differentiate under the integral, and obtain

$$(\hat{\nabla}_b)_x^{\alpha_1}(\hat{\nabla}_b)_y^{\alpha_2}T^{(3)}(x,y)| \lesssim \int_{\vartheta(x,z) > r/2} \vartheta(x,z)^{-3-|\alpha_1|} \vartheta(x,z)^{-3-|\alpha_2|} dz \lesssim r^{-2-|\alpha|}$$

This proves the desired differential inequalities for T(x, y).

Next, to prove the cancellation conditions for T, suppose ϕ is a normalized bump function in a ball $B(x_0, r)$. Then we let χ be a normalized bump function supported in $B(x_0, 4r)$, that is identically 1 on $B(x_0, 3r)$, and write

$$T\phi = T_1(\chi T_2\phi) + T_1((1-\chi)T_2\phi).$$

Now by cancellation conditions for T_2 , one sees that $r^{-1}\chi T_2\phi$ is a normalized bump function on $B(x_0, 4r)$. Hence $T_1(\chi T_2\phi)$ obeys the desired bound, namely

$$\|\hat{\nabla}^{\alpha}_{b}T_{1}(\chi T_{2}\phi)\|_{L^{\infty}(B(x_{0},r))} \lesssim_{\alpha} r^{2-|\alpha|}$$

Furthermore, for $x \in B(x_0, r)$,

$$\begin{split} \hat{\nabla}_{b}^{\alpha} T_{1}[(1-\chi)T_{2}\phi](x) \\ &= \int_{z \notin B(x,3r)} (\hat{\nabla}_{b})_{x}^{\alpha} T_{1}(x,z)(1-\chi)(z)T_{2}\phi(z)\hat{m}(z) \\ &= \int_{z \notin B(x,3r)} \int_{y \in B(x,2r)} (\hat{\nabla}_{b})_{x}^{\alpha} T_{1}(x,z)(1-\chi)(z)T_{2}(z,y)\phi(y)\hat{m}(y)\hat{m}(z). \end{split}$$

Putting absolute values,

$$\begin{aligned} |\hat{\nabla}^{\alpha}_{b}T_{1}[(1-\chi)T_{2}\phi](x)| \\ &\leq \int_{z\notin B(x,3r)} \int_{y\in B(x,2r)} \vartheta(x,z)^{-3-|\alpha|} \vartheta(z,y)^{-3} \hat{m}(y) \hat{m}(z) \\ &\lesssim \int_{z\notin B(x,3r)} \int_{y\in B(x,2r)} \vartheta(x,z)^{-6-|\alpha|} \hat{m}(y) \hat{m}(z) \\ &\lesssim r^{2-|\alpha|}, \end{aligned}$$

the second to last line following since $\vartheta(x, z) \simeq \vartheta(z, y)$ on the support of the integrals. This provides the desired bound for $\|(\hat{\nabla}_b)^{\alpha}T\phi\|_{L^{\infty}(B(x,r))}$. A similar argument establishes the bound for $\|(\hat{\nabla}_b)^{\alpha}T^*\phi\|_{L^{\infty}(B(x,r))}$. This completes our proof when both T_1 and T_2 are smoothing of order 1. Finally, suppose T_1 is smoothing of order 0, and T_2 is smoothing of order 1. Then $T := T_1 \circ T_2$ maps $C^{\infty}(\hat{X})$ continuously into itself, and so does $T^* = T_2^* \circ T_1^*$; one can repeat the above argument to show that both T and T^* satisfy the cancellation conditions for an operator of order 1. Thus it remains to compute the kernel of T, and to establish differential inequalities for the kernel of T, to which we now turn.

For $x, y \in \hat{X}$ with $x \neq y$, let $r := \vartheta(x, y)/4$, and

$$T(x,y) := (T_1k_y)(x)$$
 where $k_y(z) := T_2(z,y)$.

We first show that T(x, y) is smooth away from the diagonal, and that it satisfies the differential inequalities

$$\left| (\hat{\nabla}_b)_x^{\alpha_1} (\hat{\nabla}_b)_y^{\alpha_2} T(x, y) \right| \lesssim r^{-3 - |\alpha|}.$$

To do so, fix $x \neq y$, and let χ_1, χ_2 be normalized bump functions in B(x, r) and B(y, r) respectively, such that $\chi_1 \equiv 1$ on $B(x, r/2), \chi_2 \equiv 1$ on B(y, r/2). Then

$$T(x,y) = T_1(\chi_1 k_y)(x) + T_1(\chi_2 k_y)(x) + \int_{\hat{X}} T_1(x,z)(1-\chi_1-\chi_2)(z)T_2(z,y)\hat{m}(z).$$

The last term can be differentiated in both x and y under the integral, and the desired estimates follow. Thus it remains to consider the first two terms. But in the first term, by continuity of T on $C^{\infty}(\hat{X})$, one can differentiate with respect to y, and obtain

$$(\hat{\nabla}_b)_y^{\alpha_2}[T_1(\chi_1 k_y)(x)] = T_1[\chi_1(\hat{\nabla}_b)_y^{\alpha_2} k_y](x);$$

the latter is T_1 acting on $r^{-3-|\alpha_2|}$ times a normalized bump function in B(x,r). Thus by cancellation condition on T_1 ,

$$|(\hat{\nabla}_b)_x^{\alpha_1}(\hat{\nabla}_b)_y^{\alpha_2}[T_1(\chi_1 k_y)(x)]| \lesssim r^{-3-|\alpha|}.$$

Similarly,

$$T_1(\chi_2 k_y)(x) = \int_{\hat{X}} T_1(x, z) \chi_2(z) T_2(z, y) \hat{m}(z)$$

= $\overline{T}_2^* [T_1(x, \cdot) \chi_2(\cdot)](y).$

By continuity of \overline{T}_2^* on $C^{\infty}(\hat{X})$, one can differentiate with respect to x, and obtain

$$(\hat{\nabla}_b)_x^{\alpha_1}[T_1(\chi_2 k_y)(x)] = \overline{T}_2^*[(\hat{\nabla}_b)_x^{\alpha_1} T_1(x, \cdot)\chi_2(\cdot)](y).$$

The latter is \overline{T}_2^* acting on an $r^{-4-|\alpha_1|}$ times a normalized bump function in B(y, r), so by cancellation condition on \overline{T}_2^* , we have

$$|(\hat{\nabla}_b)_x^{\alpha_1}(\hat{\nabla}_b)_y^{\alpha_2}[T_1(\chi_2 k_y)(x)]| \lesssim r^{-3-|\alpha|}.$$

This proves our desired estimates.

It remains to show that T(x, y) is the kernel of the operator T, in the sense that (2.2) holds for all $f \in C^{\infty}(\hat{X})$ and all x not in the support of f. In fact, fix such an f, and a closed set K disjoint from the support of f. Let $\chi \in C^{\infty}(\hat{X})$ such that $\chi = 1$ on a neighborhood of K, and $\chi = 0$ on the support of f. Then for $x \in K$,

$$Tf(x) = T_1(\chi T_2 f)(x) + T_1((1-\chi)T_2 f)(x)$$

The second term is equal to

$$\int_{\hat{X}} T_1(x,z)(1-\chi)(z)T_2f(z)\hat{m}(z)$$

= $\int_{\hat{X}} \left(\int_{\hat{X}} T_1(x,z)(1-\chi)(z)T_2(z,y)\hat{m}(z) \right) f(y)\hat{m}(y),$

the last equality following from Fubini's theorem. We claim that for almost every $x \in \hat{X}$, the first term is equal to

$$\int_{\hat{X}} T_1(\chi k_y)(x) f(y) \hat{m}(y), \qquad (2.4)$$

where $k_y(z) := T_2(z, y)$; if this were true, then (2.2) holds for almost every $x \in K$. Since K is an arbitrary compact set disjoint from the support of f, (2.2) holds for almost every x not in the support of f. But then by continuity of T(x, y), and bounded convergence theorem, (2.2) holds for every x not in the support of f. Our theorem then follows.

To prove our claim, we approximate

$$\chi(z)T_2f(z) = \int_{\hat{X}} \chi(z)T_2(z,y)f(y)\hat{m}(y)$$

by Riemann sums; since T(z, y) is smooth away from the diagonal, and f is smooth, by uniform continuity, the Riemann sums converge uniformly to $\chi T_2 f$. By continuity of T in $L^2(\hat{m})$, we have T_1 of the Riemann sums converging in L^2 to $T_1(\chi T_2 f)$. Thus by passing to a subsequence, T_1 of the Riemann sums converge almost everywhere to $T_1(\chi T_2 f)$. On the other hand, T_1 of the Riemann sums is the Riemann sums of (2.4); by continuity of $T_1(\chi k_y)(x)$ for $(x, y) \in K \times \text{supp } f$, the Riemann sums of (2.4) converges uniformly to (2.4). This establishes our claim. \Box

Proof of Theorem 2.4. Suppose T is a smoothing operator of order j, $0 \le j < 4$, and $g \in \mathcal{E}(\hat{\rho}^{-\gamma}), j < \gamma < 4$. Fix $k \in \mathbb{N}_0$ and fix a point $x_0 \ne p$ sufficiently close to p. Let $r = \frac{1}{4}\vartheta(x_0, p)$, and η be a normalized bump function supported in $B(x_0, r)$, with $\eta = 1$ on $B(x_0, r/2)$. Then,

$$|\hat{\nabla}_{b}^{k}Tg(x_{0})| \leq |\hat{\nabla}_{b}^{k}T(\eta g)(x_{0})| + |\hat{\nabla}_{b}^{k}T((1-\eta)g)(x_{0})|$$

and $r^{\gamma}(\eta g)(x)$ is a normalized bump function on $B(x_0, r)$. So by the cancellation condition for T, we see that

$$|\hat{\nabla}_b^k T(\eta g)(x_0)| \le C_k r^{j-\gamma-k},$$

where $C_k > 0$ is a constant independent of x_0 and r. By using the kernel estimates, $\hat{\nabla}_b^k T((1-\eta)g)(x_0)$ can be estimated by writing out the integrals directly:

$$\hat{\nabla}_b^k T((1-\eta)g)(x_0) = \int (\hat{\nabla}_b^k T)(x_0, y)(1-\eta)(y)g(y)\hat{m}(y),$$

which can be split into two pieces. The first is over where $\vartheta(y,p) \leq r$; this piece is dominated by

$$D_k \int_{\vartheta(y,p) < r} r^{-4+j-k} \vartheta(y,p)^{-\gamma} \hat{m}(y) \le D_k r^{j-\gamma-k},$$

where $D_k > 0$ is a constant independent of x_0 and r. The second piece is over where $\vartheta(y, p) > r$; note since we have cut off those y near x_0 with $1 - \eta$ already, we can assume that $\vartheta(y, x_0) > r/2$ on this piece of integral as well. As a result, $\vartheta(x_0, y) \simeq \vartheta(y, p)$; it follows that this piece is bounded by

$$\int_{\vartheta(y,p)>r} \vartheta(y,p)^{-4+j-k} \vartheta(y,p)^{-\gamma} \hat{m}(y) \le E_k r^{j-\gamma-k}$$

where $E_k > 0$ is a constant independent of x_0 and r. Altogether,

$$|\hat{\nabla}_b^k Tg(x_0)| \le \widetilde{C}_k r^{j-\gamma-k}$$

as desired, where $\widetilde{C}_k > 0$ is a constant independent of x_0 and r. This completes our proof. \Box

We now turn to the proof of Theorem 2.2. The key is the following L^2 estimate, which can be proved by microlocalization and integration by parts (see e.g. Kohn [16]). Suppose $\hat{\overline{Z}}$ is a local section of $T^{0,1}\hat{X}$ with $\langle \hat{\overline{Z}} | \hat{\overline{Z}} \rangle_{\hat{\theta}} = 1$ on some ball $B(x, 2r) \subset \hat{X}$, and $\hat{\overline{Z}}^*$ be its formal adjoint under $L^2(\hat{m})$.

Proposition 2.5. If $\hat{\overline{Z}}^* v = u$ on B(x, 2r), where $u, v \in C^{\infty}(\hat{X})$, then for every $k \in \mathbb{N}_0$, there is a constant $C_k > 0$ independent of r and x such that

$$\begin{aligned} \|\hat{\nabla}_{b}^{k}u\|_{L^{2}(B(x,r))} \\ &\leq C_{k}(\|\hat{\nabla}_{b}^{k-1}\hat{Z}u\|_{L^{2}(B(x,2r))} + r^{-k}\|u\|_{L^{2}(B(x,2r))} + r^{-(k+1)}\|v\|_{L^{2}(B(x,2r))}). \end{aligned}$$
(2.5)

Proposition 2.6. If $\hat{\overline{Z}}v = u$ on B(x,2r), where $u, v \in C^{\infty}(\hat{X})$, then for every $k \in \mathbb{N}_0$, there is a constant $C_k > 0$ independent of r and x such that

$$\begin{aligned} \|\hat{\nabla}_{b}^{k}u\|_{L^{2}(B(x,r))} \\ &\leq C_{k}(\|\hat{\nabla}_{b}^{k-1}\hat{Z}^{*}u\|_{L^{2}(B(x,2r))} + r^{-k}\|u\|_{L^{2}(B(x,2r))} + r^{-(k+1)}\|v\|_{L^{2}(B(x,2r))}). \end{aligned}$$
(2.6)

Here the L^2 norms are taken using the norms of $L^2(\hat{m})$. Various variants and refinements of these estimates are very well-known; however, we have not been able to locate a precise reference for these estimates. For completeness and the convenience of the reader, we present the proofs of these estimates in Appendices A and B.

Using these L^2 estimates, one can prove that $\hat{\Pi}$ and \hat{K} (more precisely, the local representations \hat{K}_i 's defined by (1.36)) map $C^{\infty}(\hat{X})$ continuously into $C^{\infty}(\hat{X})$, and that their kernels satisfy differential inequalities of the correct order; cf. Christ [6,7]. In particular, if $\hat{\Pi}(x, y)$ and $\hat{K}_i(x, y)$ denote the Schwartz kernels of $\hat{\Pi}$ and \hat{K}_i respectively, then they are smooth away from the diagonal, and

$$\begin{aligned} |(\hat{\nabla}_b)_x^{\alpha_1}(\hat{\nabla}_b)_y^{\alpha_2}\hat{\Pi}(x,y)| &\lesssim_{\alpha} \vartheta(x,y)^{-4-|\alpha|}, \\ |(\hat{\nabla}_b)_x^{\alpha_1}(\hat{\nabla}_b)_y^{\alpha_2}\hat{K}_i(x,y)| &\lesssim_{\alpha} \vartheta(x,y)^{-3-|\alpha|}. \end{aligned}$$

It thus remains to prove cancellation conditions for $\hat{\Pi}$ and \hat{K}_i . The proofs will be based on strategies similar to those used in the proofs of the above kernel estimates.

Let ϕ be a normalized bump function in B(x, r). We claim that

$$\|\hat{\nabla}_{b}^{k}\hat{\Pi}\phi\|_{L^{\infty}(B(x,r))} \le C_{k}r^{-k}.$$
(2.7)

This will follow from the continuity of $\hat{\Pi}$ on $C^{\infty}(\hat{X})$ if r is sufficiently large. Therefore, without loss of generality, we assume that $r < r_0/2$, where r_0 is some small absolute constant, so that one can find a section $\hat{\overline{Z}}$ of $T^{1,0}\hat{X}$ that does not vanish on B(x,2r) (cf. discussion before (1.36)). We further normalize $\hat{\overline{Z}}$ so that $\langle \hat{\overline{Z}} | \hat{\overline{Z}} \rangle_{\hat{\theta}} = 1$ on B(x,2r). Now (2.7) is the same as showing $\|\hat{\nabla}_b^k(I - \hat{\Pi})\phi\|_{L^{\infty}(B(x,r))} \leq C_k r^{-k}$. Let v be such that $\hat{\overline{\partial}}_b^* v = (I - \hat{\Pi})\phi$, with v orthogonal to the kernel of $\hat{\overline{\partial}}_b^*$. Then by (2.5), we have

$$\begin{aligned} \|\hat{\nabla}_{b}^{k}(I-\hat{\Pi})\phi\|_{L^{2}(B(x,r))} \\ &\leq C_{k}(\|\hat{\nabla}_{b}^{k-1}\hat{\overline{Z}}\phi\|_{L^{2}(B(x,2r))} + r^{-k}\|(I-\hat{\Pi})\phi\|_{L^{2}(B(x,2r))} + r^{-(k+1)}\|\langle v|\hat{\omega}\rangle_{\hat{\theta}}\|_{L^{2}(B(x,2r))}), \end{aligned}$$

where $\hat{\omega}$ is the dual (0,1) form to $\hat{\overline{Z}}$ on B(x,2r). The first term on the right hand side is bounded by

$$C_k r^{-k} |B(x,2r)|^{1/2} = C_k r^{2-k}.$$

In the second term, we estimate $||(I - \hat{\Pi})\phi||_{L^2(B(x,2r))}$ trivially by $||\phi||_{L^2(\hat{X})} \leq C|B(x,r)|^{1/2} = Cr^2$ (recall ϕ is normalized in B(x,r)), so that the second term is bounded by $C_k r^{2-k}$ as well. In the last term, we estimate using Poincaré-type inequality (see Corollary 11.5* of Christ [6]):

$$\|\langle v|\hat{\omega}\rangle_{\hat{\theta}}\|_{L^{2}(B(x,2r))} \leq Cr\|(I-\hat{\Pi})\phi\|_{L^{2}(\hat{X})} \leq Cr\|\phi\|_{L^{2}(\hat{X})} \leq Crr^{2}.$$

Thus altogether,

$$\|\hat{\nabla}_{b}^{k}(I-\hat{\Pi})\phi\|_{L^{2}(B(x,r))} \leq C_{k}r^{2-k},$$

and since this holds for all k, by Sobolev embedding,

$$\|\hat{\nabla}_b^k(I-\hat{\Pi})\phi\|_{L^{\infty}(B(x,r))} \le C_k r^{-k}$$

as desired. (2.7) follows, and since $\hat{\Pi}$ is self-adjoint on $L^2(\hat{m})$, this completes the proof that $\hat{\Pi}$ is smoothing of order 0.

Let now $\hat{\Pi}_1: L^2_{(0,1)}(\hat{m}, \hat{\theta}) \to L^2_{(0,1)}(\hat{m}, \hat{\theta})$ be the Szegö projection on (0,1) forms, i.e. the orthogonal projection onto the kernel of $\hat{\overline{\partial}}_b^*$ in $L^2_{(0,1)}(\hat{m}, \hat{\theta})$. Using the partition of unity given just before (1.36), we can define local representations of $\hat{\Pi}_1$, by letting

$$(\hat{\Pi}_1)_{ij}\varphi := \langle \eta_i \hat{\Pi}_1(\eta_j \varphi \hat{\omega}_j) | \hat{\omega}_i \rangle_{\hat{\theta}}.$$

Then $(\hat{\Pi}_1)_{ij}$ sends functions to functions for all $1 \leq i, j \leq N$, and

$$\hat{\Pi}_{1}\phi = \sum_{i,j} \left((\hat{\Pi}_{1})_{ij} [\langle \phi | \hat{\omega}_{j} \rangle_{\hat{\theta}}] \right) \hat{\omega}_{i}$$

for any (0, 1) form ϕ on \hat{X} . A proof similar to the above shows that $\hat{\Pi}_1$ is a smoothing operator of order 0, in the sense that the local representations $(\hat{\Pi}_1)_{ij}$ are all smoothing of order 0; for instance, to prove that $(\hat{\Pi}_1)_{ij}$ satisfies the desired cancellation conditions, if φ is a normalized bump function on a sufficiently small ball B(x, 2r) that intersects the support of η_i , one would apply Proposition 2.6 with \hat{Z} being the (0, 1) vector field dual to $\hat{\omega}_i$, $u = \eta_j \varphi \langle \hat{\omega}_j | \hat{\omega}_i \rangle_{\hat{\theta}} - \langle \Pi_1(\eta_j \varphi \hat{\omega}_j) | \hat{\omega}_i \rangle_{\hat{\theta}}$, and v being a function that solves $\hat{\bar{\partial}}_b v = (I - \hat{\Pi}_1)(\eta_j \varphi \hat{\omega}_j)$. In fact then $\hat{Z}v = u$ on B(x, 2r). We omit the details.

Now, let φ be a normalized bump function in a ball B(x, r) that intersects the support of η_i , with $r < r_0/4$ as before. We prove cancellation properties for \hat{K}_i and \hat{K}_i^* , namely

$$\|\hat{\nabla}_b^k \hat{K}_i \varphi\|_{L^{\infty}(B(x,r))} \le C_k r^{1-k}, \qquad (2.8)$$

and

$$\|\hat{\nabla}_b^k \hat{K}_i^* \varphi\|_{L^{\infty}(B(x,r))} \le C_k r^{1-k}.$$
(2.9)

To prove the former, let

$$u = (I - \hat{\Pi})(\psi \hat{K}_i \varphi),$$

where $\psi \equiv 1$ on B(x, 2r), and ψ is a normalized bump function on B(x, 4r). We apply estimate (2.5) for this *u*. On B(x, 2r), $\hat{\bar{\partial}}_{b}u = (I - \hat{\Pi}_{1})(\eta_{i}\varphi\hat{\omega}_{i})$, since ψ is identically 1 there. In other words, writing \hat{Z} for the dual of $\hat{\omega}_{i}$, we have

$$\hat{\overline{Z}}u = \langle (I - \hat{\Pi}_1)(\eta_i \varphi \hat{\omega}_i) | \hat{\omega}_i \rangle_{\hat{\theta}} = \eta_i \varphi - \sum_j [(\hat{\Pi}_1)_{ji} \varphi] \langle \hat{\omega}_j | \hat{\omega}_i \rangle_{\hat{\theta}}$$

on B(x,2r). So the first term on the right hand side of (2.5) is bounded by $Cr^{1-k}|B(x,2r)|^{1/2} = Cr^{1-k}r^2$, by the cancellation property of $\hat{\Pi}_1$ we just proved above. On the other hand, by Proposition B of Christ [7], since u is orthogonal to the kernel of $\hat{\overline{\partial}}_b$,

$$||u||_{L^2(B(x,2r))} \le Cr ||\hat{\overline{\partial}}_b u||_{L^2_{(0,1)}(\hat{X})},$$

which implies

$$\|u\|_{L^{2}(B(x,2r))} \leq Cr\left(\|\psi(I-\hat{\Pi}_{1})(\eta_{i}\varphi\hat{\omega}_{i})\|_{L^{2}_{(0,1)}(\hat{X})} + \|(\hat{\nabla}_{b}\psi)\hat{K}_{i}\varphi\|_{L^{2}(\hat{X})}\right).$$

But the first term in the bracket is bounded by $\|\eta_i\varphi\|_{L^2(\hat{X})} \leq C|B(x,r)|^{1/2} = Cr^2$, and the second term is bounded by Cr^2 by the kernel estimates on \hat{K}_i (note φ is supported on B(x,r), while $\hat{\nabla}_b \psi$ is supported in an annulus $B(x,4r) \setminus B(x,2r)$). This in turn implies

$$|u||_{L^2(B(x,2r))} \le Crr^2,$$

and the second term on the right hand side of (2.5) is bounded by $Cr^{1-k}r^2$. Finally, let v be such that $\hat{\overline{\partial}}_b^* v = u$, so that $\hat{\overline{Z}}^*[\langle v|\hat{\omega}_i\rangle_{\hat{\theta}}] = u$ on B(x, 2r). Then

$$\|\langle v|\hat{\omega}_i\rangle_{\hat{\theta}}\|_{L^2(B(x,2r))} \le Cr\|u\|_{L^2(\hat{X})} \le Cr\|\psi\hat{K}_i\varphi\|_{L^2(\hat{X})} \le Cr^2\|\varphi\|_{L^2(\hat{X})} \le Cr^2r^2.$$

(The first and the fourth inequality are both applications of Proposition B of Christ [7] again.) Thus altogether, $\|\hat{\nabla}_{b}^{k}u\|_{L^{2}(B(x,r))} \leq C_{k}r^{1-k}r^{2}$ for all k, and by Sobolev embedding, this implies

$$\|\hat{\nabla}_b^k u\|_{L^\infty(B(x,r))} \le C_k r^{1-k}$$

for all k.

Remember we want the same estimate for $\hat{K}_i \varphi$ in place of u, so as to prove (2.8). But $\hat{K}_i \varphi - u$ can be computed on B(x, r) fairly easily. In fact, since $\hat{K}_i \varphi = \hat{K}(\eta_i \varphi \hat{\omega}_i)$ is orthogonal to the kernel of $\overline{\partial}_b$, we have $\hat{K}_i \varphi = (I - \hat{\Pi}) \hat{K}_i \varphi$. So

$$\hat{K}_i \varphi - u = (I - \hat{\Pi})(1 - \psi)\hat{K}_i \varphi.$$

Since $\psi \equiv 1$ on B(x, 2r), we have

$$\hat{K}_i \varphi - u = -\hat{\Pi}(1 - \psi)\hat{K}_i \varphi$$
 on $B(x, r)$

It follows that for $y \in B(x, r)$,

$$\hat{\nabla}_b^k(\hat{K}_i\varphi - u)(y) = -\sum_{j=1}^\infty \int_{2^j r \le \vartheta(z,x) \le 2^{j+1}r} (\hat{\nabla}_b)_y^k \hat{\Pi}(y,z)(1-\psi)(z)\hat{K}_i\varphi(z)\hat{m}(z),$$

 \mathbf{SO}

$$\begin{aligned} \|\hat{\nabla}_{b}^{k}(\hat{K}_{i}\varphi-u)\|_{L^{\infty}(B(x,r))} &\leq C_{k}\sum_{j=1}^{\infty}(2^{j}r)^{-2-k}\|\hat{K}_{i}\varphi\|_{L^{2}(B(x,2^{j+1}r))} \\ &\leq C_{k}\sum_{j=1}^{\infty}(2^{j}r)^{-2-k}(2^{j}r)\|\hat{\overline{\partial}}_{b}\hat{K}_{i}\varphi\|_{L^{2}(\hat{X})}. \end{aligned}$$

(The last inequality is Proposition B of Christ [7].) By estimating the term $\|\hat{\partial}_b \hat{K}_i \varphi\|_{L^2(\hat{X})}$ by $\|(I - \hat{\Pi}_1)(\eta_i \varphi \hat{\omega}_i)\|_{L^2(\hat{X})} \leq \|\varphi\|_{L^2(\hat{X})} \leq Cr^2$, we get

$$\|\hat{\nabla}_{b}^{k}(\hat{K}_{i}\varphi - u)\|_{L^{\infty}(B(x,r))} \le C_{k}r^{-2-k}rr^{2} = C_{k}r^{1-k}$$

By combining with the previous estimate on $\hat{\nabla}_{b}^{k} u$, we get

$$\|\hat{\nabla}_b^k \hat{K}_i \varphi\|_{L^{\infty}(B(x,r))} \le C_k r^{1-k}$$

as desired. (2.8) follows then. A similar argument proves (2.9), since \hat{K}^* is the partial inverse of $\hat{\overline{\partial}}_b^*$. This shows that \hat{K} and \hat{K}^* are smoothing of order 1, and it follows now that $\hat{N} = \hat{K}\hat{K}^*$ (more precisely, $\hat{N} = \sum_{i,j} \hat{K}_i \langle \hat{\omega}_i | \hat{\omega}_j \rangle_{\hat{\theta}} \hat{K}_j^*$) is smoothing of order 2. This completes the proof of Theorem 2.2.

We conclude this section by making the following useful observation: as was demonstrated in Folland and Stein [10] (see Theorem 15.15 there), on \hat{X} one can construct some smoothing operators T_0 , T_1 of orders 1, such that schematically,

$$I = T_1 \hat{\nabla}_b + T_0.$$

Thus if A is a smoothing operator of order 1, then for any positive integers j, one can then find smoothing operators A_0, A_1, \ldots, A_j of orders 1 such that

$$\hat{\nabla}_{b}^{j}A = \sum_{i=0}^{j} A_{i}\hat{\nabla}_{b}^{i}; \qquad (2.10)$$

in fact, e.g. when j = 1, one just needs to observe

$$\hat{\nabla}_b A = \hat{\nabla}_b A (T_1 \hat{\nabla}_b + T_0),$$

and the desired equality follows by letting $A_i = \hat{\nabla}_b A T_i$, i = 0, 1. (A_i is smoothing of order 1 by Theorem 2.3 above.) The general case for (2.10) follows by induction on j. (2.10) can be thought of as a way of commuting derivatives past smoothing operators. In particular, if $NL^{k,p}$ denotes the non-isotropic Sobolev space, given by the set of all functions whose $\hat{\nabla}_b^j$ is in $L^p(\hat{m})$ for $j = 0, 1, 2, \ldots, k$, then (2.10) implies the first part of the following proposition:

Proposition 2.7.

- (a) Any smoothing operator of order 1 maps $NL^{k,p}$ continuously into $NL^{k+1,p}$, for all $k \ge 0$ and all 1 .
- (b) Any smoothing operator of order 1 maps $L^{\infty}(\hat{m})$ continuously into $L^{\infty}(\hat{m})$.

The last part of this proposition then follows from the case k = 0 of the first part by noting that $L^{\infty}(\hat{m})$ embeds into $L^{p}(\hat{m})$ for any p > 4, and that $NL^{1,p}$ embeds into $L^{\infty}(\hat{m})$ by Sobolev embedding.

3. A key mapping property

3.1. The main theorem

In this section, we prove the following theorem, which allows one to establish the important mapping property (1.25) of $(I + \hat{R}^{*,\widetilde{m}}\chi)^{-1}$. We use the notion of smoothing operators of order j we introduced in the last section.

Theorem 3.1. Suppose A is a smoothing operator of order 1 on \hat{X} , and h is a function in $\mathcal{E}(\hat{\rho}^1)$ supported in a sufficiently small neighborhood of p. We write h also for the operator

that is multiplication by h. Then the bounded linear operator $I - Ah: L^2(\hat{m}) \to L^2(\hat{m})$ is invertible, and its inverse extends to a continuous linear map

$$(I - Ah)^{-1}$$
: $\mathcal{E}(\hat{\rho}^{-4+\delta}) \to \mathcal{E}(\hat{\rho}^{-4+\delta})$

for every $0 < \delta < 4$.

The key of the proof is the following:

Lemma 3.2. Suppose A and h are as in the above theorem. Then for any non-negative integer k, and any function $u \in C^{\infty}(X) \cap L^{\infty}(\hat{X})$, we have

$$\|\hat{\rho}^k \hat{\nabla}_b^k (Ah)^{k+1} u\|_{L^{\infty}(\hat{m})} \lesssim_k \|u\|_{L^{\infty}(\hat{m})}.$$

Assuming the lemma, we first prove Theorem 3.1.

Proof of Theorem 3.1. First, if A is smoothing of order 1, then since \hat{X} is compact, the kernel of A satisfies

$$\sup_{x\in\hat{X}}\int_{\hat{X}}|A(x,y)|\hat{m}(y)+\sup_{y\in\hat{X}}\int_{\hat{X}}|A(x,y)|\hat{m}(x)<\infty.$$

It follows that A is bounded on $L^2(\hat{m})$. If the support of h is a sufficiently small neighborhood of p, then since $h \in \mathcal{E}(\hat{\rho}^1)$, one can make $||h||_{L^{\infty}}$ sufficiently small. Thus the norm of Ah, as a bounded linear operator on $L^2(\hat{m})$, can be made smaller than 1/2. This in turn allows one to invert I - Ah by a Neumann series: for $u \in L^2(\hat{m})$, one has

$$u + (Ah)u + (Ah)^{2}u + (Ah)^{3}u + \dots$$

converging to a limit v in $L^2(\hat{m})$, and (I - Ah)v = u. Thus (I - Ah) is invertible on $L^2(\hat{m})$, and its inverse is given by the Neumann series

$$(I - Ah)^{-1} = I + Ah + (Ah)^{2} + (Ah)^{3} + \dots$$

Now we extend $(I - Ah)^{-1}$ to $\mathcal{E}(\hat{\rho}^{-4+\delta})$, $0 < \delta < 4$. In order to do so, we need to further assume that the norm of Ah, as a bounded linear operator on $L^{\infty}(\hat{m})$, is smaller than 1/2. That can be achieved if the support of h is sufficiently small.

Suppose now $u \in \mathcal{E}(\hat{\rho}^{-4+\delta}), 0 < \delta < 4$. Let $v = [I + (Ah) + (Ah)^2 + ...]u$. We want to show that $v \in \mathcal{E}(\hat{\rho}^{-4+\delta})$. To do so, suppose k is a non-negative integer. To show that

$$\hat{\rho}^{k+(4-\delta)}\hat{\nabla}^k_b v \in L^{\infty}(\hat{m}), \tag{3.1}$$
we split the sum defining v into two parts: let $v_1 = [I + (Ah) + \cdots + (Ah)^{k+2}]u$, and $v_2 = [(Ah)^{k+3} + (Ah)^{k+4} + \ldots]u$. Then $\hat{\rho}^{k+(4-\delta)}\hat{\nabla}_b^k v_1 \in L^{\infty}(\hat{m})$ by Theorem 2.4, since each term of v_1 is in $\mathcal{E}(\hat{\rho}^{-4+\delta})$. Furthermore, note that

$$A: \mathcal{E}(\hat{\rho}^{-1+\delta}) \to L^{\infty}(\hat{m}) \quad \text{for all } 0 < \delta < 1.$$
(3.2)

This holds because $\mathcal{E}(\hat{\rho}^{-1+\delta}) \subset L^p(\hat{m})$ for some p > 4, and $A: L^p(\hat{m}) \to L^\infty(\hat{m})$ whenever p > 4. Thus from $u \in \mathcal{E}(\hat{\rho}^{-4+\delta}), 0 < \delta < 4$, we conclude, from Theorem 2.4 and (3.2), that $(Ah)^2 u \in C^\infty(X) \cap L^\infty(\hat{X})$. As a result, by Lemma 3.2, and our bound of Ah on $L^\infty(\hat{m})$, we have

$$\begin{aligned} \|\hat{\rho}^k \hat{\nabla}_b^k v_2\|_{L^{\infty}(\hat{m})} &\leq C_k \sum_{\ell=0}^{\infty} \|(Ah)^{2+\ell} u\|_{L^{\infty}(\hat{m})} \\ &\leq C_k \sum_{\ell=0}^{\infty} 2^{-\ell} \|(Ah)^2 u\|_{L^{\infty}(\hat{m})} \leq C_k \end{aligned}$$

Combining this with the bound for v_1 , (3.1) follows, and this shows $v \in \mathcal{E}(\hat{\rho}^{-4+\delta})$ as desired. \Box

3.2. An auxiliary family of operators

Now we need to detour into a discussion of a two-parameter family of operators, that will be indexed by two non-negative integers j and ℓ . Suppose A is a smoothing operator of order 1. Suppose also that $q_{\ell}(x, y)$ is a function in $C^{\infty}(\hat{X} \times \hat{X})$ that vanishes to non-isotropic order ℓ along the diagonal, i.e.

$$|q_{\ell}(x,y)| \lesssim \vartheta(x,y)^{\ell}$$

for some non-negative integer ℓ . We will write $q_{\ell}^{x}(y) := q_{\ell}(x, y)$; by abuse of notation, we will also denote by q_{ℓ}^{x} the multiplication operator $v(y) \mapsto q_{\ell}^{x}(y)v(y)$. Given a non-negative integer j, for $v \in C^{\infty}(\hat{X})$ and $x \in \hat{X}$, we define

$$Tv(x) = \left. \left(\hat{\nabla}_b \right)_z^j \right|_{z=x} [Aq_\ell^x v](z).$$
(3.3)

This is well defined, since $q_{\ell}^{x}v$ is a C^{∞} function on \hat{X} for each fixed x, and A maps $C^{\infty}(\hat{X})$ into $C^{\infty}(\hat{X})$. We will see below that this assignment $v \mapsto Tv$ defines a continuous map from $C^{\infty}(\hat{X})$ to $C^{\infty}(\hat{X})$. Since the properties of this map depend mainly only on the integers j and ℓ , we will denote any operator of this form by $T_{j,\ell}$. In other words, if $v \in C^{\infty}(\hat{X})$ and $x \in \hat{X}$, then $T_{j,\ell}v(x)$ is given by the right hand side of (3.3) for some smoothing operator A of order 1, and some $q_{\ell} \in C^{\infty}(\hat{X} \times \hat{X})$ that vanishes to non-isotropic order ℓ along the diagonal. **Lemma 3.3.** Suppose $v \in C^{\infty}(\hat{X})$, and for all $x \in \hat{X}$, we have

$$T_{j,\ell}v(x) := \left. (\hat{\nabla}_b)_z^j \right|_{z=x} [Aq_\ell^x v](z) \tag{3.4}$$

for some A and q_{ℓ} as above. Then $T_{j,\ell}v$ is a C^{∞} function on \hat{X} . Furthermore, for any $r \geq 0$, one has the following Leibniz rule:

$$(\hat{\nabla}_b)_x^r T_{j,\ell} v(x) = \sum_{s=0}^r \binom{r}{s} \left(\hat{\nabla}_b \right)_z^{j+r-s} \Big|_{z=x} \left[A(\hat{\nabla}_b)_x^s q_\ell^x v \right](z)$$
(3.5)

where $(\hat{\nabla}_b)_x^s q_\ell^x$ denotes the multiplication operator $v(y) \mapsto (\hat{\nabla}_b)_x^s q_\ell(x, y)v(y)$. In addition, the adjoint $T_{j,\ell}^*$ of $T_{j,\ell}$ with respect to $L^2(\hat{m})$ maps $C^{\infty}(\hat{X})$ into itself, and is given by

$$T_{j,\ell}^* w(z) = A^* [(\hat{\nabla}_b^*)^j (\bar{q}_{\ell,z} w)](z)$$
(3.6)

where $q_{\ell,z}(x) := q_{\ell}(x,z)$.

Proof. Suppose $v \in C^{\infty}(\hat{X})$, and $x \in \hat{X}$. First we show that $T_{j,\ell}$ is differentiable at x, and that (3.5) holds when r = 1. In fact, for any smooth curve $\gamma: (-1, 1) \to \hat{X}$ with $\gamma(0) = x, \gamma'(0) = Y$, we have

$$\frac{1}{\varepsilon} \left[T_{j,\ell} v(\gamma(\varepsilon)) - T_{j,\ell} v(\gamma(0)) \right] \\
= \left. \left(\hat{\nabla}_b \right)_z^j \Big|_{z=\gamma(\varepsilon)} \left[A \left(\frac{q_\ell^{\gamma(\varepsilon)} v - q_\ell^{\gamma(0)} v}{\varepsilon} - Y_x q_\ell^x \right) v \right] (z) \\
+ \frac{1}{\varepsilon} \left(\left(\hat{\nabla}_b \right)_z^j \Big|_{z=\gamma(\varepsilon)} - \left(\hat{\nabla}_b \right)_z^j \Big|_{z=\gamma(0)} \right) [A q_\ell^x v] (z) \\
+ \left. \left(\hat{\nabla}_b \right)_z^j \Big|_{z=\gamma(\varepsilon)} [A (Y_x q_\ell^x) v] (z).$$
(3.7)

(We wrote Y_x to emphasize that the derivative is with respect to x.) Now the first term on the right hand side of (3.7) is bounded by

$$\left\| (\hat{\nabla}_b)_z^j A\left(\frac{q_\ell^{\gamma(\varepsilon)} v - q_\ell^{\gamma(0)}}{\varepsilon} - Y_x q_\ell^x \right) v(z) \right\|_{L^{\infty}(\hat{m}(z))},$$

which tends to zero as $\varepsilon \to 0$ since

$$\left(\frac{q_{\ell}^{\gamma(\varepsilon)}v - q_{\ell}^{\gamma(0)}}{\varepsilon} - Y_x q_{\ell}^x\right) v(z) \to 0 \quad \text{in } C^{\infty}(\hat{X}) \text{ as a function of } z,$$

and $A: C^{\infty}(\hat{X}) \to C^{\infty}(\hat{X})$ is continuous. Next, the second term on the right hand side of (3.7) converges to

$$Y_z \left(\hat{\nabla}_b \right)_z^j [Aq_\ell^x v](z) \Big|_{z=x}$$

as $\varepsilon \to 0$, since $Aq_{\ell}^x v(z)$ is C^{∞} as a function of $z \in \hat{X}$. Finally, the last term on the right hand side of (3.7) converges to

$$(\hat{\nabla}_b)_z^j\Big|_{z=x} \left[A(Y_x q_\ell^x)v\right](z)$$

as $\varepsilon \to 0$, since $(\hat{\nabla}_b)_z^j [A(Y_x q_\ell^x) v](z)$ is a continuous function of $z \in \hat{X}$. This proves $T_{j,\ell} v$ is differentiable at x, and that

$$[YT_{j,\ell}v](x) = \left(Y_z(\hat{\nabla}_b)_z^j\right)\Big|_{z=x} \left[Aq_\ell^x v\right](z) + \left(\hat{\nabla}_b\right)_z^j\Big|_{z=x} \left[A(Y_x q_\ell^x)v\right](z).$$
(3.8)

In particular, (3.5) holds when r = 1.

By successive differentiation of (3.8), using the case r = 1 of (3.5), then shows that $T_{j,\ell}v \in C^{\infty}(\hat{X})$, and that (3.5) holds for all $r \geq 1$.

Now we prove that $T_{j,\ell}^*$ is given by the expression (3.6). The case when j = 0 is easy; in fact then

$$T_{0,\ell}v(x) = \int_{\hat{X}} A(x,y)q_{\ell}(x,y)v(y)d\hat{m}(y),$$

so integrating this against $\overline{w(x)}d\hat{m}(x)$, and applying Fubini's theorem, we see that $T^*_{0,\ell}w(y) = A^*[\bar{q}_{\ell,y}w](y)$, as desired.

The case j > 0 is more tricky, because while formally

$$T_{j,\ell}v(x) = \int_{\hat{X}} [(\nabla_b)_x^j A(x,y)] q_\ell(x,y) v(y) d\hat{m}(y),$$

the integral on the right hand side simply does not converge, due to the singularity of $(\nabla_b)_x^j A(x, y)$. Hence we have to proceed differently, as follows.

We assume that (3.6) holds for all $T_{k,\ell}^*$ with $0 \le k \le j-1$. Take r = 1 and replace j to j-1 in (3.5), we have

$$T_{j,\ell} = \hat{\nabla}_b T_{j-1,\ell} - T_{j-1,\ell_0}, \qquad (3.9)$$

where $\ell_0 = \min \{ \ell - 1, 0 \}$ and

$$T_{j-1,\ell_0}v(x) := \left. (\hat{\nabla}_b)_z^{j-1} \right|_{z=x} \left[A(\hat{\nabla}_b)_x q_\ell^x v \right](z).$$

By taking adjoint of (3.9) in the sense of distribution with respect to $L^2(\hat{m})$, we deduce

$$T_{j,\ell}^* = T_{j-1,\ell}^* (\hat{\nabla}_b)^* - T_{j-1,\ell_0}^*.$$
(3.10)

From (3.10) and the induction assumptions, we can check that

$$T_{j,\ell}^* w(z) = A^* [(\hat{\nabla}_b^*)^{j-1} (\bar{q}_{\ell,z} \hat{\nabla}_b^* w)](z) - A^* [(\hat{\nabla}_b^*)^{j-1} (\hat{\nabla}_b q_{\ell,z}) w)](z)$$

= $A^* [(\hat{\nabla}_b^*)^j (\bar{q}_{\ell,z} w)](z).$

(3.6) follows. As a result, by repeating the proof of (3.5), $T_{j,\ell}^*$ maps $C^{\infty}(\hat{X})$ into $C^{\infty}(\hat{X})$, with

$$(\nabla_b)_z^r T_{j,\ell}^* w(z) = \sum_{s=0}^r \binom{r}{s} \left. (\hat{\nabla}_b)_y^{r-s} \right|_{y=z} \left[A(\hat{\nabla}_b)_z^s (\hat{\nabla}_b^*)^j (\bar{q}_{\ell,z} w) \right](y).$$

Lemma 3.4. For any $j, \ell \geq 0$, the linear operators

$$T_{j,\ell}: C^{\infty}(\hat{X}) \to C^{\infty}(\hat{X})$$

and

$$T_{j,\ell}^*: C^\infty(\hat{X}) \to C^\infty(\hat{X})$$

considered in the previous lemma are both continuous.

Proof. Suppose v_m converges to v in $C^{\infty}(\hat{X})$ as $m \to \infty$. Then for any $k, s \ge 0$,

$$\|(\hat{\nabla}_b)_z^k[(\hat{\nabla}_b)_z^s q_\ell^x(z)(v_m-v)(z)]\|_{L^2(\hat{m}(z))} \to 0 \quad \text{uniformly in } x \in \hat{X},$$

and by continuity of A on $NL^{k,2}$ (see Proposition 2.7(a)), we see that for any $k, s \ge 0$,

$$\|(\hat{\nabla}_b)_z^k[A(\hat{\nabla}_b)_x^s q_\ell^x(v_m - v)](z)\|_{L^2(\hat{m}(z))} \to 0 \quad \text{uniformly in } x \in \hat{X}.$$

By Sobolev embedding, it follows that the same is true for all $k, s \ge 0$ if the L^2 norm is replaced by L^{∞} in the above equation. Hence by (3.5), we conclude that

$$\|(\nabla_b)^r T_{j,\ell}(v_m - v)\|_{L^{\infty}(\hat{X})} \to 0$$

for all $r \geq 0$. Since this is true for all r, we proved $T_{j,\ell}v_m \to T_{j,\ell}v$ in $C^{\infty}(\hat{X})$ as $m \to \infty$. A similar argument, based on (3.6) instead, proves that $T_{j,\ell}^*: C^{\infty}(\hat{X}) \to C^{\infty}(\hat{X})$ is continuous. \Box

Lemma 3.5. For any $\ell \geq 0$, the operator $T_{\ell+1,\ell}$ is smoothing of order 0.

Proof. By the previous lemma, both $T_{\ell+1,\ell}$ and its adjoint $T^*_{\ell+1,\ell}$ map $C^{\infty}(\hat{X})$ into $C^{\infty}(\hat{X})$ continuously.

Given $v \in C^{\infty}(\hat{X})$, and $x \in \hat{X}$ not in the support of v, if z is in a sufficiently small neighborhood of x, we have z not in the support of $q_{\ell}^{x}v$. Hence if A(x, y) is the kernel of A, then for all such z, we have

$$[Aq_{\ell}^{x}v](z) = \int_{\hat{X}} A(z,y)q_{\ell}(x,y)v(y)\hat{m}(y).$$

It follows that one can differentiate under the integral, and obtain

$$T_{j,\ell}v(x) = \int_{\hat{X}} [(\hat{\nabla}_b)_x^j A(x,y)] q_\ell(x,y)v(y)\hat{m}(y).$$

Hence the kernel of $T_{j,\ell}$ is given by

$$T_{j,\ell}(x,y) = [(\hat{\nabla}_b)_x^j A(x,y)] q_\ell(x,y)$$

When $j = \ell + 1$, this kernel satisfies the differential inequalities

$$|(\hat{\nabla}_b)_x^{\alpha_1}(\hat{\nabla}_b)_y^{\alpha_2}T_{j,\ell}(x,y)| \lesssim \vartheta(x,y)^{-4-|\alpha|},$$

since A is a smoothing operator of order 1, and q_{ℓ} vanishes to order ℓ along the diagonal. It follows that the kernel of $T_{\ell+1,\ell}$ satisfies the differential inequalities of a smoothing operator of order 0.

Finally, we verify the cancellation property of $T_{\ell+1,\ell}$. To do so, suppose ϕ is a normalized bump function in some ball $B(x_0, r_0) \subset \hat{X}$. Then for any $x \in B(x_0, r_0)$ and $s \ge 0$, we have $(\hat{\nabla}_b)_x^s q_\ell^x \phi$ being $r_0^{\ell-s}$ times a normalized bump function in $B(x_0, r_0)$. Thus by cancellation property of A, we have

$$\|(\hat{\nabla}_b)^{j+r-s}[A(\hat{\nabla}_b)^s_x q^x_\ell \phi]\|_{L^{\infty}(B(x_0,r))} \lesssim r_0^{1-(j+r-\ell)}$$

for all j, r and s with $j \ge 0$, $r \ge s$. In particular, evaluating at $x \in B(x_0, r_0)$, we have

$$\left| \left(\hat{\nabla}_b \right)_y^{j+r-s} \right|_{z=x} \left[A(\hat{\nabla}_b)_x^s q_\ell^x \phi \right](z) \right| \lesssim r_0^{1-(j+r-\ell)}.$$

Hence by (3.5), we see that

$$|(\hat{\nabla}_b)_x^r T_{j,\ell} \phi(x)| \lesssim r_0^{1-(j+r-\ell)}$$

Since this is true for all $x \in B(x_0, r_0)$, we obtain

$$\|(\hat{\nabla}_b)_x^r T_{j,\ell}\phi\|_{L^{\infty}(B(x_0,r_0))} \lesssim r_0^{1-(j+r-\ell)}.$$

When $j = \ell + 1$, this gives the desired cancellation property of order 0 for $T_{\ell+1,\ell}$. Similarly, one can prove the desired cancellation property of $T^*_{\ell+1,\ell}$. Hence $T_{\ell+1,\ell}$ is a smoothing operator of order 0. \Box

Lemma 3.6. Suppose $v \in C^{\infty}(X) \cap L^{\infty}(\hat{X})$. Then for $x \in X$, one can still define $T_{j,\ell}v(x)$ by (3.4), and formula (3.5) continues to hold for all $x \in X$.

Proof. Suppose $v \in C^{\infty}(X) \cap L^{\infty}(\hat{X})$, and $x \in X$. Let $\eta \in C_0^{\infty}(X)$ be such that $\eta \equiv 1$ in a neighborhood of x, and write $w_1 = \eta v$, $w_2 = (1 - \eta)v$. Then $w_1 \in C^{\infty}(\hat{X})$, and $T_{j,\ell}w_1(x)$ is given by (3.4) with v replaced by w_1 . Furthermore, w_2 is identically zero near x. Thus by pseudolocality of A, we have $[Aq_{\ell}^x w_2](z)$ being C^{∞} for z near x, and one can define $T_{j,\ell}w_2(x)$ by (3.4) with v replaced by w_2 . Since $T_{j,\ell}v = T_{j,\ell}w_1 + T_{j,\ell}w_2$, it follows that one can define $T_{j,\ell}v(x)$ by (3.4).

In order to differentiate $T_{j,\ell}v$ at x, it suffices to differentiate $T_{j,\ell}w_1$ and $T_{j,\ell}w_2$ at x. One can differentiate $T_{j,\ell}w_1$ using (3.5). To differentiate $T_{j,\ell}w_2$ at x, note that since x is not in the support of w_2 , for any z in a small neighborhood of x, we have

$$T_{j,\ell}w_2(z) = \int_{\hat{X}} [(\hat{\nabla}_b)_z^j A(z,y)] q_\ell(z,y) w_2(y) \hat{m}(y).$$

Differentiating under the integral using the dominated convergence theorem, one sees that the derivatives of $T_{i,\ell}w_2$ at x satisfies

$$(\hat{\nabla}_b)_x^r T_{j,\ell} w_2(x) = \sum_{s=0}^r \binom{r}{s} (\hat{\nabla}_b)_z^{j+r-s} \Big|_{z=x} \left[A(\hat{\nabla}_b)_x^s q_\ell^x w_2 \right](z).$$

Together, one concludes that the derivatives of $T_{j,\ell}v(x)$ is given by (3.5) in our current case as well. \Box

3.3. Proof of Lemma 3.2

We now move on to the proof of Lemma 3.2. To do so, we fix a neighborhood U of p, and fix a frame \hat{Z}_1 , $\hat{Z}_{\bar{1}}$, T of the complexified tangent bundle in U as in Section 1.5. Write X_1 and X_2 for the real and imaginary parts of \hat{Z}_1 . One can now define normal coordinates centered at any point $x \in U$: for y sufficiently close to x, there exists a unique $w \in \mathbb{R}^3$ such that if $\gamma(t)$ is the integral curve of $w_1X_1 + w_2X_2 + w_3T$ with $\gamma(0) = x$, then $\gamma(1) = y$. In that case we write $y = x \exp(w)$, or equivalently $w = \Theta(x, y)$. Note in this case,

$$\vartheta(x,y) \simeq |w_1| + |w_2| + |w_3|^{1/2}.$$

Also, in this normal coordinate system, X_1 and X_2 then takes the form

$$X_{1} = \frac{\partial}{\partial w_{1}} + 2w_{2}\frac{\partial}{\partial w_{3}} + O^{1}\frac{\partial}{\partial w_{1}} + O^{1}\frac{\partial}{\partial w_{2}} + O^{2}\frac{\partial}{\partial w_{3}},$$
$$X_{2} = \frac{\partial}{\partial w_{2}} - 2w_{1}\frac{\partial}{\partial w_{3}} + O^{1}\frac{\partial}{\partial w_{1}} + O^{1}\frac{\partial}{\partial w_{2}} + O^{2}\frac{\partial}{\partial w_{3}}$$

where O^k are functions that vanish to non-isotropic order $\geq k$ at w = 0. It follows that if α is the multi-index $(\alpha_1, \alpha_2, \alpha_3)$ and $j \leq ||\alpha||$ (here $||\alpha|| := |\alpha_1| + |\alpha_2| + 2|\alpha_3|$ is the non-isotropic length of α), then $(\hat{\nabla}_b)_x^j \Theta(x, y)^{\alpha}$ vanishes to non-isotropic order $||\alpha|| - j$ along the diagonal y = x, i.e.

$$|(\hat{\nabla}_b)_x^j \Theta(x,y)^{\alpha}| \lesssim \vartheta(x,y)^{\|\alpha\|-j}.$$

Now fix $h \in \mathcal{E}(\hat{\rho}^1)$. For $x \in U$, let h^x be the function of w defined by $h^x(w) = h(x \exp(w))$. Let P_k^x be the Taylor polynomial of h^x at w = 0 up to non-isotropic order k. We sometimes think of P_k^x as a function of y. Then

$$P_k^x(y) = \sum_{\|\alpha\| \le k} \frac{1}{\alpha!} (\partial_w^{\alpha} h^x)(0) [\Theta(x,y)]^{\alpha}.$$

One can show, by reduction to the ordinary Taylor's theorem, that

$$|h(y) - P_{k-1}^{x}(y)| \le C\hat{\rho}(x)^{1-k}\vartheta(x,y)^{k} \quad \text{if } \vartheta(x,y) < \frac{1}{4}\hat{\rho}(x).$$

$$(3.11)$$

We will make crucial use of this in the proof of the following lemma:

Lemma 3.7. Suppose A is a smoothing operator of order 1, and $h \in \mathcal{E}(\hat{\rho}^1)$. If $v \in C^{\infty}(X) \cap L^{\infty}(\hat{X})$, then for any $k \geq 1$, we have

$$\left\| \hat{\rho}^{k}(x) \left(\hat{\nabla}_{b} \right)_{z}^{k} \right\|_{z=x} \left[A(h - P_{k-1}^{x}) v \right](z) \right\|_{L^{\infty}(\hat{m}(x))} \lesssim \|v\|_{L^{\infty}(\hat{m})}$$

where $P_{k-1}^{x}(y)$ is defined as above.

Proof. Fix $x \in X$, and write $r := \hat{\rho}(x)/8$. Let $\varepsilon \in (0, r)$, and write

$$1 = (1 - \varphi_r) + (\varphi_r - \varphi_\varepsilon) + \varphi_\varepsilon$$

where φ_r and φ_{ε} are normalized bump functions on B(x, 2r) and $B(x, 2\varepsilon)$ respectively, with $\varphi_r(y) \equiv 1$ on B(x, r), $\varphi_{\varepsilon}(y) \equiv 1$ on $B(x, \varepsilon)$. Then

$$h - P_{k-1}^{x} = (h - P_{k-1}^{x})(1 - \varphi_{r}) + (h - P_{k-1}^{x})(\varphi_{r} - \varphi_{\varepsilon}) + (h - P_{k-1}^{x})\varphi_{\varepsilon},$$

and we estimate the contribution of each of these three terms to

$$\hat{\rho}^{k}(x) \left. (\hat{\nabla}_{b})_{z}^{k} \right|_{z=x} [A(h - P_{k-1}^{x})v](z).$$

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Let's call the above contributions I, II and III. We will show that I and II are bounded by $C||v||_{L^{\infty}(\hat{m})}$ uniformly in x and ε , while III tends to 0 as $\varepsilon \to 0$. These obviously imply the desired conclusion in our lemma.

Now note I and II can be computed using the kernel of A: we have

$$I = \hat{\rho}^{k}(x) \int_{\hat{X}} (1 - \varphi_{r}(y)) [(\hat{\nabla}_{b})_{x}^{k} A(x, y)] [h(y) - P_{x}^{k-1}(y)] v(y) \hat{m}(y)$$

and

$$II = \hat{\rho}^k(x) \int_{\hat{X}} (\varphi_r(y) - \varphi_\varepsilon(y)) [(\hat{\nabla}_b)_x^k A(x,y)] [h(y) - P_x^{k-1}(y)] v(y) \hat{m}(y).$$

I can then be estimated by breaking up $h - P_{k-1}^x$ into h and the individual terms in P_{k-1}^x : the term involving h only is bounded by

$$\begin{split} \hat{\rho}^{k}(x) & \int_{\vartheta(y,x)\geq r} |(\hat{\nabla}_{b})_{x}^{k}A(x,y)h(y)v(y)|\hat{m}(y) \\ &\leq \hat{\rho}^{k}(x) \int_{\vartheta(y,x)\geq r} \vartheta(x,y)^{-3-k}\hat{m}(y) \|v\|_{L^{\infty}(\hat{m})} \\ &\lesssim \begin{cases} r\|v\|_{L^{\infty}(\hat{m})} & \text{if } k > 1 \\ r\log r\|v\|_{L^{\infty}(\hat{m})} & \text{if } k = 1 \\ &\lesssim \|v\|_{L^{\infty}(\hat{m})}. \end{split}$$

Also, the term involving P_{k-1}^x can be bounded by

$$\begin{split} \hat{\rho}^k(x) & \int\limits_{\vartheta(y,x) \ge r} |[(\hat{\nabla}_b)^k_x A(x,y)] P^x_{k-1}(y) v(y) | \hat{m}(y) \\ & \leq \sum_{\|\alpha\| \le k-1} \hat{\rho}^k(x) \int\limits_{\vartheta(y,x) \ge r} \vartheta(x,y)^{-3-k} \hat{\rho}^{1-\|\alpha\|}(x) \vartheta(x,y)^{\|\alpha\|} \hat{m}(y) \|v\|_{L^{\infty}(\hat{m})}. \end{split}$$

The terms when $\|\alpha\| = k - 1$ can be bounded by

$$\hat{\rho}^2(x) \int_{\vartheta(y,x) \ge r} \vartheta(x,y)^{-4} \hat{m}(y) \|v\|_{L^{\infty}(\hat{m})} \lesssim r^2 \log r \|v\|_{L^{\infty}(\hat{m})} \lesssim \|v\|_{L^{\infty}(\hat{m})}.$$

The terms when $\|\alpha\| < k - 1$ can be bounded by

$$\hat{\rho}(x)^{k+1-\|\alpha\|} \int_{\vartheta(x,y)\geq r} \vartheta(x,y)^{-4+\|\alpha\|-(k-1)} \hat{m}(y) \|v\|_{L^{\infty}(\hat{m})}$$
$$\lesssim \hat{\rho}(x)^{k+1-\|\alpha\|} r^{\|\alpha\|-(k-1)} \|v\|_{L^{\infty}(\hat{m})}$$
$$\lesssim r^{2} \|v\|_{L^{\infty}(\hat{m})} \lesssim \|v\|_{L^{\infty}(\hat{m})}.$$

This shows that $|I| \lesssim ||v||_{L^{\infty}(\hat{m})}$ uniformly in x, as desired.

Next, II can be bounded using (3.11), uniformly in x and ε :

$$\begin{split} \hat{\rho}^{k}(x) & \int_{\vartheta(y,x) < r} |[(\hat{\nabla}_{b})^{k}_{x}A(x,y)](h - P^{x}_{k-1})(y)v(y)|\hat{m}(y) \\ & \leq \hat{\rho}^{k}(x) \int_{\vartheta(y,x) < r} \vartheta(x,y)^{-3-k}r^{1-k}\vartheta(x,y)^{k}\hat{m}(y)\|v\|_{L^{\infty}(\hat{m})} \\ & \lesssim r^{2}\|v\|_{L^{\infty}(\hat{m})} \lesssim \|v\|_{L^{\infty}(\hat{m})}. \end{split}$$

Thus it remains to show that III tends to 0 as $\varepsilon \to 0$.

To do so, note that there is some constant C_v (possibly depending on many derivatives of v) such that $C_v^{-1}\varepsilon^{-k}r^k\varphi_{\varepsilon}(h-P_{k-1}^x)v$ is a normalized bump function in $B(x, 2\varepsilon)$. In fact, the $L^{\infty}(\hat{m})$ norm of $\varepsilon^{-k}r^k\varphi_{\varepsilon}(h-P_{k-1}^x)$ is bounded by

$$C\varepsilon^{-k}r^kr^{1-k}\varepsilon^k \le C.$$

Also, each $\hat{\nabla}_b$ derivative of this function gains at worst ε^{-1} , and the function v is C^{∞} on this ball $B(x, 2\varepsilon)$. Hence $C_v^{-1}\varepsilon^{-k}r^k\varphi_{\varepsilon}(h-P_{k-1}^x)v$ is a normalized bump function in $B(x, 2\varepsilon)$, as claimed. Since A is smoothing of order 1, by the cancellation conditions of A, we have

$$|III| \le r^k \left\| (\hat{\nabla}_b)^k [A\varphi_{\varepsilon}(h - P_{k-1}^x)v] \right\|_{L^{\infty}(B(x, 2\varepsilon))} \le C_v \varepsilon,$$

which tends to 0 as $\varepsilon \to 0$. This completes the proof of the current lemma. \Box

We are now ready to prove Lemma 3.2. In fact we will prove the following slight generalization:

Lemma 3.8. For any non-negative integer k, if $A_1, A_2, \ldots, A_{k+1}$ are smoothing operators of order 1, and $h \in \mathcal{E}(\hat{\rho}^1)$ is supported in a sufficiently small neighborhood of p, then for any function $u \in C^{\infty}(X) \cap L^{\infty}(\hat{X})$, we have

$$\|\hat{\rho}^k \nabla_b^k (A_{k+1}hA_kh\dots A_1h)u\|_{L^{\infty}(\hat{m})} \lesssim_k \|u\|_{L^{\infty}(\hat{m})}.$$

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For simplicity, we will write below S_k for the operator $A_kh \ldots A_1h$. We then need to bound $\|\hat{\rho}^k \hat{\nabla}_b^k S_{k+1}u\|_{L^{\infty}(\hat{m})}$ for $u \in C^{\infty}(X) \cap L^{\infty}(\hat{X})$.

Proof. We proceed by induction on k. When k = 0, one just need to notice that $||A_1hu||_{L^{\infty}(\hat{m})} \leq ||u||_{L^{\infty}(\hat{m})}$, which holds since A_1 and h both preserve $L^{\infty}(\hat{m})$. Suppose now the proposition has been proved up to k-1 for some positive integer k. In other words, we assume

$$\|\hat{\rho}^{i}\hat{\nabla}_{b}^{i}S_{i+1}u\|_{L^{\infty}(\hat{m})} \lesssim_{i} \|u\|_{L^{\infty}(\hat{m})}$$
(3.12)

for all $0 \leq i \leq k - 1$. Let $u \in C^{\infty}(X) \cap L^{\infty}(\hat{X})$, and write

$$\hat{\rho}^k \hat{\nabla}^k_b S_{k+1} u(x) = \hat{\rho}^k \hat{\nabla}^k_b A_{k+1} h S_k u(x) = V_1(x) + V_2(x)$$

where

$$V_1(x) = \hat{\rho}^k(x) \left. \left(\hat{\nabla}_b \right)_z^k \right|_{z=x} \left[A_{k+1} (h - P_{k-1}^x) S_k u \right](z)$$

and

$$V_2(x) = \hat{\rho}^k(x) \left. \left(\hat{\nabla}_b \right)_z^k \right|_{z=x} [A_{k+1} P_{k-1}^x S_k u](z).$$

We want to show that both $V_1(x)$ and $V_2(x)$ are bounded by $C||u||_{L^{\infty}(\hat{m})}$.

First $V_1(x)$ is like an error term, that can be estimated by Lemma 3.7. In fact, applying it to $v = S_k u$, which is in $C^{\infty}(X) \cap L^{\infty}(\hat{X})$, we have

$$\left\| \hat{\rho}^k(x) \left(\hat{\nabla}_b \right)_z^k \right\|_{z=x} \left[A_{k+1} (h - P_{k-1}^x) (S_k u) \right](z) \right\|_{L^{\infty}(\hat{m}(x))} \le C \| S_k u \|_{L^{\infty}(\hat{m})}.$$

and the latter is bounded trivially by $C ||u||_{L^{\infty}(\hat{m})}$. Thus it remains to estimate the main term $V_2(x)$.

Now write $\Theta_x(y) = \Theta(x, y)$. Then

$$V_2(x) = \sum_{\|\alpha\| \le k-1} \frac{1}{\alpha!} (\partial_w^{\alpha} h^x)(0) \hat{\rho}^k(x) \left(\hat{\nabla}_b \right)_z^k \Big|_{z=x} \left[A_{k+1} \Theta_x^{\alpha} S_k u \right](z).$$

Note that on the right hand side of this sum, $(\hat{\nabla}_b)_y^k \Big|_{z=x} [A_{k+1}\Theta_x^{\alpha}S_k u](z)$ is of the form $T_{k,\ell}S_k u$ if $\ell = \|\alpha\|$, where $T_{k,\ell}$ is defined as in (3.4), with $A = A_{k+1}$ and $q_\ell(x,y) = \Theta(x,y)^{\alpha}$. Thus

$$|V_2(x)| \lesssim \sum_{\ell=0}^{k-1} \hat{\rho}^{k-\ell+1}(x) |T_{k,\ell}(S_k u)(x)|.$$

We will prove, by induction, that when $1 \leq j \leq k$, and $0 \leq \ell \leq j - 1$, if $T_{j,\ell}$ is any operator of the form (3.4), we have

$$\left\|\hat{\rho}^{j-\ell-1}(x)T_{j,\ell}(S_k u)(x)\right\|_{L^{\infty}(\hat{m}(x))} \lesssim \|u\|_{L^{\infty}(\hat{m})},\tag{3.13}$$

and that the bound is uniform in the choice of A and q_{ℓ} , as long as A is normalized as a smoothing operator of order 1 (so that all constants arising in the definition of a smoothing operator of order 1 are uniformly bounded), and q_{ℓ} is normalized as a smooth function vanishing to non-isotropic order ℓ along the diagonal (so that all constants arising in the definition of such functions are uniformly bounded). Assume this for the moment. Then applying it to the case j = k, with $T_{k,\ell}$ defined using $A = A_{k+1}$ and $q_{\ell}(x,y) = \Theta(x,y)^{\alpha}$, we have

$$|V_2(x)| \lesssim \sum_{\ell=0}^{k-1} \hat{\rho}^2(x) ||u||_{L^{\infty}(\hat{m})} \lesssim ||u||_{L^{\infty}(\hat{m})}.$$

This will complete the proof of the current lemma.

It remains to prove (3.13). To do so, we proceed by induction on j. First, when j = 1 (hence $\ell = 0$), note that $T_{1,0}$ is a smoothing operator of order 0 by Lemma 3.5. Hence

$$T_{1,0}S_ku = (T_{1,0}A_k)hS_{k-1}u = \tilde{A}hS_{k-1} = \tilde{S}_1S_{k-1}u,$$

where $\tilde{A} := T_{1,0}A_k$ is smoothing of order 1 by Theorem 2.3, and $\tilde{S}_1 = \tilde{A}h$. Since \tilde{S}_1 and S_{k-1} both preserve $L^{\infty}(\hat{m})$ (which can be see, e.g. by our induction hypothesis (3.12) with i = 0, and its iteration), we see that

$$||T_{1,0}S_ku(x)||_{L^{\infty}(\hat{m}(x))} = ||S_1[S_{k-1}u](x)||_{L^{\infty}(\hat{m}(x))} \lesssim ||S_{k-1}u||_{L^{\infty}(\hat{m})} \lesssim ||u||_{L^{\infty}(\hat{m})}$$

Hence (3.13) holds when j = 1.

Next, fix positive integer j_0 with $2 \le j_0 \le k$, and assume that (3.13) has been proved for all $j \le j_0 - 1$. In other words, we assume that we have verified already

$$\left\|\hat{\rho}^{a-b-1}(x)T_{a,b}(S_k u)(x)\right\|_{L^{\infty}(\hat{m}(x))} \lesssim \|u\|_{L^{\infty}(\hat{m})}$$
(3.14)

for all a and b with $1 \leq a \leq j_0 - 1$ and $0 \leq b \leq a - 1$. We want to prove the same statement when $a = j_0$. So we fix ℓ with $0 \leq \ell \leq j_0 - 1$. First, $S_k u \in C^{\infty}(X) \cap L^{\infty}(\hat{X})$ since u is as such. Thus we can apply Lemma 3.6 and (3.5) to $v = S_k u$: in fact, if $T_{j_0,\ell}$ is defined so that $T_{j_0,\ell}v(x) = (\hat{\nabla}_b)_z^{j_0}\Big|_{z=x} [Aq_\ell^x v](z)$, then by (3.5), one can compute $\hat{\nabla}_b^{j_0-\ell-1}[T_{\ell+1,\ell}(S_k u)](x)$ as a sum over s from 0 to $j_0 - \ell - 1$. In this sum, the term with s = 0 is exactly $T_{j_0,\ell}(S_k u)(x)$. Hence by rearranging, we have

$$T_{j_0,\ell}(S_k u)(x) = \hat{\nabla}_b^{j_0-\ell-1}[T_{\ell+1,\ell}(S_k u)](x) - \sum_{s=1}^{j_0-\ell-1} {j_0-\ell-1 \choose s} (\hat{\nabla}_b)_z^{j_0-s} \Big|_{z=x} [A(\hat{\nabla}_b)_x^s q_\ell^x S_k u](z). \quad (3.15)$$

The first term on the right hand side can then be written as

$$\hat{\nabla}_{b}^{j_{0}-\ell-1}[\tilde{A}hS_{k-1}u](x) = \hat{\nabla}_{b}^{j_{0}-\ell-1}[\tilde{A}hS_{j_{0}-\ell-1}](S_{k-j_{0}+\ell}u)(x)$$

where $\tilde{A} := T_{\ell+1,\ell}A_k$ is smoothing of order 1 by Lemma 3.5 and Theorem 2.3. Hence by our induction hypothesis (3.12), with $i = j_0 - \ell - 1$ (note $0 \le i \le k - 1$ under our assumptions so (3.12) applies), we have

$$|\hat{\rho}^{j_0-\ell-1}(x)\hat{\nabla}_b^{j_0-\ell-1}[T_{\ell+1,\ell}(S_k u)](x)| \lesssim \|S_{k-j_0+\ell}u\|_{L^{\infty}(\hat{m})} \lesssim \|u\|_{L^{\infty}(\hat{m})}.$$
 (3.16)

To bound the rest of the sum on the right hand side of (3.15), note that $(\hat{\nabla}_b)_z^{j_0-s}\Big|_{z=x} \times [A(\hat{\nabla}_b)_x^s q_\ell^x S_k u](z)$ is of the form $T_{j_0-s,\max\{\ell-s,0\}}S_k u(x)$. This is because $(\hat{\nabla}_b)_x^s q_\ell(x,y)$ is smooth on $\hat{X} \times \hat{X}$, and vanishes on the diagonal to order $\max\{\ell-s,0\}$. Hence

$$\left| \hat{\rho}^{j_{0}-\ell-1}(x) \sum_{s=1}^{j_{0}-\ell-1} \left(\frac{j_{0}-\ell-1}{s} \right) (\hat{\nabla}_{b})_{z}^{j_{0}-s} \Big|_{z=x} \left[A(\hat{\nabla}_{b})_{x}^{s} q_{\ell}^{x} S_{k} u \right](z) \right|$$

$$\lesssim \sum_{s=1}^{j_{0}-\ell-1} \hat{\rho}^{\max\{s-\ell,0\}}(x) \hat{\rho}^{j_{0}-s-\max\{\ell-s,0\}-1}(x) |T_{j_{0}-s,\max\{\ell-s,0\}} S_{k} u(x)|$$

$$\lesssim \sum_{s=1}^{j_{0}-\ell-1} \hat{\rho}^{\max\{s-\ell,0\}}(x) ||u||_{L^{\infty}(\hat{m})} \lesssim ||u||_{L^{\infty}(\hat{m})}, \qquad (3.17)$$

the second inequality following from our induction hypothesis (3.14) with $a = j_0 - s$, $b = \max\{\ell - s, 0\}$ (note that then $1 \le a \le j_0 - 1$ and $0 \le b \le j_0 - s - 1$ in the above sum over s, so (3.14) applies for this a and b). Combining (3.15), (3.16) and (3.17), we see that

$$|\rho^{j_0 - \ell - 1}(x) T_{j_0, \ell} S_k u(x)| \lesssim ||u||_{L^{\infty}(\hat{m})}$$

uniformly in x, as desired. This completes our proof of (3.13), and hence the proof of the current lemma. \Box

4. Construction of the CR function ψ

The goal in this section is to show that there exists a function $\psi \in C^{\infty}(\hat{X})$ such that $\hat{\overline{\partial}}_b \psi = 0$ and near p, we have $\psi = |z|^2 + it + R$, $R = \varepsilon(\hat{\rho}^4)$ (see Theorem 4.4). It is via

this ψ that we reduce our problem from the non-compact manifold X to the compact manifold \hat{X} , as was explained in the introduction.

Until further notice, we work in some small neighborhood of p. Put

$$\overline{Z}_1^0 = \frac{\partial}{\partial \overline{z}} + iz \frac{\partial}{\partial t}.$$

We say that g is a quasi-homogeneous polynomial of degree $d \in \mathbb{N}_0 := \mathbb{N} \bigcup \{0\}$ if g is the finite sum

$$g = \sum_{\alpha+\beta+2\gamma=d,\alpha,\beta,\gamma\in\mathbb{N}_0} c_{\alpha,\beta,\gamma} \overline{z}^{\alpha} z^{\beta} t^{\gamma}, \quad c_{\alpha,\beta,\gamma}\in\mathbb{C}.$$

First, we need

Lemma 4.1. For any monomial $a\overline{z}^{\alpha}z^{\beta}t^{\gamma}$, $a \in \mathbb{C}$, $\alpha, \beta, \gamma \in \mathbb{N}_0$, with $\alpha + \beta + 2\gamma \geq 3$, we can find a polynomial f, that is quasi-homogeneous of degree $\alpha + \beta + 2\gamma + 1$, such that $\overline{Z}_1^0 f = az^{\alpha}\overline{z}^{\beta}t^{\gamma}$ and

$$|\operatorname{Re} f(z,t)| \le c |z|^2 (|z|+|t|) \text{ in a neighborhood of } p,$$

where $c > 0$ is a constant. (4.1)

Proof. We proceed by induction over γ . First we assume that $\gamma = 0$. Given a monomial $a\overline{z}^{\alpha}z^{\beta}$, $a \in \mathbb{C}$, $\alpha, \beta \in \mathbb{N}_0$, with $\alpha + \beta \geq 3$. Put $f = \frac{a}{\alpha+1}\overline{z}^{\alpha+1}z^{\beta}$. It is easy to see that $\overline{Z}_1^0 f = a\overline{z}^{\alpha}z^{\beta}$ and (4.1) hold. Let $\gamma \geq 1$. Given a monomial $a\overline{z}^{\alpha}z^{\beta}t^{\gamma}$, $a \in \mathbb{C}$, $\alpha, \beta, \gamma \in \mathbb{N}_0$, with $\alpha + \beta + 2\gamma \geq 3$. First, we assume that $\alpha = \beta = 0$. By the induction assumption, we can find a quasi-homogeneous polynomial f_1 , of degree $2\gamma + 1$, such that

$$\overline{Z}_1^0 f_1 = -ia\gamma \left| z \right|^2 t^{\gamma - 1} + i\overline{a}\gamma z^2 t^{\gamma - 1}$$

and (4.1) hold. Put

$$f = a\overline{z}t^{\gamma} - \overline{a}zt^{\gamma} + f_1.$$

It is not difficult to check that $\overline{Z}_1^0 f = at^{\gamma}$ and (4.1) hold. Now, we assume that $\alpha + \beta \geq 1$. By the induction assumption, we can find a quasi-homogeneous polynomial f_1 , of degree $\alpha + \beta + 2\gamma + 1$, such that

$$\overline{Z}_1^0 f_1 = -i \frac{a\gamma}{\alpha+1} \overline{z}^{\alpha+1} z^{\beta+1} t^{\gamma-1}$$

and (4.1) hold. Put

$$f = \frac{a}{\alpha+1}\overline{z}^{\alpha+1}z^{\beta}t^{\gamma} + f_1.$$

It is not difficult to check that $\overline{Z}_1^0 f = a \overline{z}^{\alpha} z^{\beta} t^{\gamma}$ and (4.1) hold.

The lemma follows. $\hfill \square$

Proposition 4.2. There exists a function $\varphi \in C^{\infty}(\hat{X})$ such that $\operatorname{Re} \varphi \geq 0$ on \hat{X} , $\overline{\partial}_b \varphi$ vanishes to infinite order at p and near p, we have

$$\varphi(z,t) = (2\pi)(|z|^2 + it) + S, \quad S = \varepsilon(\hat{\rho}^6).$$

Proof. We assume that the local coordinates (z, t) defined on a small open set $W \subset X$ of p. From Lemma 4.1 and (1.31), it is not difficult to see that we can find $f_j(z,t)$, $j = 6, 7, \ldots$, where for each j, f_j is a quasi-homogeneous polynomial of degree j, such that

$$\hat{Z}_{\bar{1}}((2\pi)(|z|^2 + it) + \sum_{j=6}^m f_j(z,t)) \in \varepsilon(\hat{\rho}^{m+3}), \quad m = 6, 7, \dots,$$
(4.2)

and for each j = 6, 7, ...,

$$|\operatorname{Re} f_j(z,t)| \le c_j |z|^2 (|z|+|t|) \text{ on } W_j \subset W, \ c_j > 0 \text{ is a constant},$$

$$(4.3)$$

where W_j is an open set, for each $j = 6, 7, \ldots$ Take $\phi(z, t) \in C_0^{\infty}(\mathbb{C} \times \mathbb{R}, \mathbb{R}_+)$ so that $\phi(z, t) = 1$ if $|z^2| + |t| \leq \frac{1}{2}$ and $\phi(z, t) = 0$ if $|z^2| + |t| \geq 1$. For each $j = 6, 7, \ldots$, take $\epsilon_j > 0$ be a small constant ($\epsilon_j \sim 2^{-j}$ will do) so that $\operatorname{Supp} \phi(\frac{z}{\epsilon_i}, \frac{t}{\epsilon_i}) \subset W_j$,

$$\left|\phi(\frac{z}{\epsilon_j}, \frac{t}{\epsilon_j^2}) \operatorname{Re} f_j(z, t)\right| < 2^{-j} |z|^2$$
(4.4)

and for all $\alpha, \beta, \gamma \in \mathbb{N}_0$, $\alpha + \beta + 2\gamma < j$, we have

$$\left\| \partial_{\bar{z}}^{\alpha} \partial_{z}^{\beta} \partial_{t}^{\gamma} \left(\phi(\frac{z}{\epsilon_{j}}, \frac{t}{\epsilon_{j}^{2}}) f_{j}(z, t) \right) \right\|_{L^{\infty}} < 2^{-j}.$$

$$(4.5)$$

On W, we put

$$\varphi_1(z,t) = (2\pi)(|z|^2 + it) + \sum_{j=6}^{\infty} \phi(\frac{z}{\epsilon_j}, \frac{t}{\epsilon_j^2}) f_j(z,t).$$

From (4.5), we can check that $\varphi_1(z,t)$ is well-defined as a smooth function on W and for all $\alpha, \beta, \gamma \in \mathbb{N}_0$, $\alpha + \beta + 2\gamma = d$, $d \ge 6$, we have

$$\partial_{\bar{z}}^{\alpha}\partial_{z}^{\beta}\partial_{t}^{\gamma}\varphi_{1}|_{(0,0)} = \partial_{\bar{z}}^{\alpha}\partial_{z}^{\beta}\partial_{t}^{\gamma}f_{d}|_{(0,0)}.$$

Combining this with (4.2), we conclude that $\hat{\bar{\partial}}_b \varphi_1$ vanishes to infinite order at p. Moreover, from (4.4), we have

Re
$$\varphi_1(z,t) \ge |z^2| (2\pi - \sum_{j=6}^{\infty} 2^{-j}) > \frac{1}{2} |z|^2$$

Thus, $\operatorname{Re} \varphi_1 \geq 0$ on W. Take $\chi \in C_0^{\infty}(W)$, $\chi \geq 0$ and $\chi = 1$ near p and put $\varphi = \chi \varphi_1 \in C^{\infty}(X)$. Then, φ satisfies the claim of this proposition. The proposition follows. \Box

Let W be a small neighborhood of p such that $|\varphi(z,t)| \ge \hat{\rho}(z,t)^2$ on W, where φ is as in Proposition 4.2. Take $\chi \in C_0^{\infty}(W, \mathbb{R}_+)$ so that $\chi = 1$ in some small neighborhood of p. Put

$$\tau := \frac{\chi}{\varphi}.$$

It is easy to check that τ is well-defined as an element in $\mathscr{D}'(\hat{X})$. Since $\hat{\bar{\partial}}_b \varphi$ vanishes to infinite order at p, we have $\hat{\bar{\partial}}_b \tau \in C^{\infty}(\hat{X})$. Thus,

$$\hat{\Box}_b \tau \in C^\infty(\hat{X}).$$

Now

$$I = \hat{\Pi} + \hat{N}\hat{\Box}_b \quad \text{on } \mathcal{D}'(\hat{X}), \tag{4.6}$$

since $I = \hat{\Pi} + \hat{\Box}_b \hat{N}$ on $L^2(\hat{m})$, and since the operators involved are self-adjoint and pseudolocal. Hence we obtain

$$\tau = \hat{\Pi}\tau + \hat{N}\hat{\Box}_b\tau = \hat{\Pi}\tau - F, \qquad (4.7)$$

where $F \in C^{\infty}(\hat{X})$. Thus,

$$\overline{\partial}_b(\tau + F) = 0 \quad \text{in } \mathscr{D}'(\hat{X}). \tag{4.8}$$

Take $C_0 > 0$ be a large constant so that $\operatorname{Re} F + C_0 > 0$ on \hat{X} . Since $\operatorname{Re} \varphi \ge 0$, we conclude that

$$\operatorname{Re}(\tau + F + C_0) > 0 \text{ on } \hat{X}.$$
 (4.9)

Put

$$\psi = \frac{1}{\tau + F + C_0}.$$

From (4.9), we know that $\tau + F + C_0 \neq 0$ on X. Since $\tau + F + C_0 \in C^{\infty}(X)$, we conclude that $\psi \in C^{\infty}(X)$. Now, we study the behavior of ψ near p. Let $W' \Subset W$ be a small neighborhood of p such that $\chi = 1$ on W' and $|\varphi(F + C_0)| < 1$ on W'. Then, on W',

$$\psi = \frac{1}{\tau + F + C_0} = \frac{1}{\frac{1}{\varphi} + F + C_0} = \frac{\varphi}{1 + \varphi(F + C_0)} \in C^{\infty}(W').$$
(4.10)

Thus,

$$\psi \in C^{\infty}(\hat{X}). \tag{4.11}$$

Moreover, from (4.10), we can check that near p,

$$\psi = (2\pi)(|z|^2 + it) + R, \quad R \in \varepsilon(\hat{\rho}^4).$$
 (4.12)

Lemma 4.3. We have

$$\hat{\overline{\partial}}_b \psi = 0 \quad on \ \hat{X}$$

Proof. Put $h := \tau + F + C_0$, so that $\psi = h^{-1}$, and take any $g \in \Omega_0^{0,1}(X)$. We have

$$(\hat{\overline{\partial}}_b \psi \mid g)_{\hat{m},\hat{\theta}} = -(\hat{\overline{\partial}}_b h \mid \overline{h}^{-2}g)_{\hat{m},\hat{\theta}} = 0$$
(4.13)

since $\hat{\bar{\partial}}_b h = 0$ in the sense of distribution. Thus, $\hat{\bar{\partial}}_b \psi = 0$ on X. Since $\psi \in C^{\infty}(\hat{X})$, we conclude that $\hat{\bar{\partial}}_b \psi = 0$ on \hat{X} . The lemma follows. \Box

From (4.9), (4.11), (4.12) and Lemma 4.3, we obtain the main result of this section.

Theorem 4.4. There is a smooth function $\psi \in C^{\infty}(\hat{X})$ such that $\overline{\partial}_b \psi = 0$ on \hat{X} , $\psi \neq 0$ on X, Re $\psi \geq 0$ on \hat{X} and near p, we have

$$\psi(z,t) = (2\pi)(|z|^2 + it) + R, \quad R = \varepsilon(\hat{\rho}^4).$$
 (4.14)

In the study of the positive *p*-mass theorem (see [5]), one needs to find a special CR function of specific growth rate on \hat{X} . More precisely, in [5], one needs to find a CR function $g \in C^{\infty}(X)$ with

$$g = \frac{z}{\left|z\right|^2 + it} + g_1$$

near p, where $g_1 \in \mathcal{E}(\hat{\rho}^0)$. By using the proof of Theorem 4.4, we can construct a such CR function:

Theorem 4.5. There is a function $g \in C^{\infty}(X) \bigcap \mathscr{D}'(\hat{X})$ such that $\hat{\overline{\partial}}_b g = 0$ and

$$g = \frac{1}{\psi}(z+r),$$

where ψ is as in Theorem 4.4 and $r \in C^{\infty}(\hat{X})$, $r = \varepsilon(\hat{\rho}^2)$.

Proof. We can repeat the proof of Proposition 4.2 with minor change and conclude that there is a function $\tilde{r} \in C^{\infty}(\hat{X})$ with $\tilde{r} = \varepsilon(\hat{\rho}^5)$ such that $\hat{\overline{\partial}}_b(z+\tilde{r})$ vanishes to infinite order at p. Put

$$\widetilde{g} = \frac{1}{\psi}(z + \widetilde{r}).$$

Since $\hat{\bar{\partial}}_b(z+\tilde{r})$ vanishes to infinite order at p and $\hat{\bar{\partial}}_b\psi=0$, we have

$$\bar{\bar{\partial}}_b \tilde{g} \in C^\infty(\hat{X}). \tag{4.15}$$

Thus,

$$\hat{\Box}_b \widetilde{g} \in C^\infty(\widehat{X}).$$

By (4.6), we obtain

$$\widetilde{g} = \widehat{\Pi}\widetilde{g} + \widehat{N}\widehat{\Box}_b\widetilde{g} = \widehat{\Pi}\widetilde{g} - f,$$

where $f \in C^{\infty}(\hat{X})$. Thus,

$$\hat{\overline{\partial}}_b(\widetilde{g}+f) = 0 \text{ on } \mathscr{D}'(\hat{X}).$$

Put $g = \tilde{g} + f$. Then $\hat{\bar{\partial}}_b g = 0$ and we can check that $g = \frac{1}{\psi}(z + \tilde{r} + \psi f) = \frac{1}{\psi}(z + r)$, where $r = \tilde{r} + \psi f \in C^{\infty}(\hat{X}), r = \varepsilon(\hat{\rho}^2)$. The theorem follows. \Box

5. The relation between $\Box_{b,1}$ and $\widetilde{\Box}_b$

Having constructed our CR function ψ , we can proceed as in Section 1.4, and construct the intermediate Kohn Laplacian $\widetilde{\Box}_b$. We refer the reader to that section for the details of this construction. The goal of the current section is to reduce the study of our operator of interest, namely $\Box_{b,1}$, to the study of $\widetilde{\Box}_b$. This is done by establishing (1.6) and (1.7) in Section 1.4.

First, we recall that $m_1 = |\psi|^{-2} \widetilde{m}$. Thus,

$$u \in L^{2}(m_{1}) \text{ if and only if } \psi^{-1}u \in L^{2}(\widetilde{m}),$$

$$u \in L^{2}_{(0,1)}(m_{1},\hat{\theta}) \text{ if and only if } \psi^{-1}u \in L^{2}_{(0,1)}(\widetilde{m},\hat{\theta}),$$
(5.1)

and

$$\begin{aligned} \|u\|_{m_{1}} &= \left\|\psi^{-1}u\right\|_{\widetilde{m}}, \quad \|v\|_{m_{1},\hat{\theta}} = \left\|\psi^{-1}v\right\|_{\widetilde{m},\hat{\theta}}, \forall u \in L^{2}(m_{1}), v \in L^{2}_{(0,1)}(m_{1},\hat{\theta}), \\ (u_{1} \mid u_{2})_{m_{1}} &= (\psi^{-1}u_{1} \mid \psi^{-1}u_{2})_{\widetilde{m}}, \quad \forall u_{1}, u_{2} \in L^{2}(m_{1}), \\ (v_{1} \mid v_{2})_{m_{1},\hat{\theta}} &= (\psi^{-1}v_{1} \mid \psi^{-1}v_{2})_{\widetilde{m},\hat{\theta}}, \quad \forall v_{1}, v_{2} \in L^{2}_{(0,1)}(m_{1},\hat{\theta}). \end{aligned}$$
(5.2)

Let

$$\overline{\partial}_{b,1}^{*,f}:\Omega^{0,1}(X)\to C^{\infty}(X) \tag{5.3}$$

be the formal adjoint of $\overline{\partial}_b$ with respect to $(\cdot | \cdot)_{m_1}, (\cdot | \cdot)_{m_1,\hat{\theta}}$. Let also

$$\widetilde{\overline{\partial}}_b^{*,f}: \Omega^{0,1}(X) \to C^\infty(X) \tag{5.4}$$

be the formal adjoint of $\overline{\partial}_b$ with respect to $(\cdot | \cdot)_{\widetilde{m}}, (\cdot | \cdot)_{\widetilde{m},\hat{\theta}}$. Then

Lemma 5.1. For $v \in \Omega_0^{0,1}(X)$, we have

$$\overline{\partial}_{b,1}^{*,f}v = \psi \widetilde{\overline{\partial}}_b^{*,f}(\psi^{-1}v).$$
(5.5)

Proof. Let $h \in C_0^{\infty}(X)$, $g \in \Omega_0^{0,1}(X)$. We have

$$(\overline{\partial}_b h | g)_{m_1,\hat{\theta}} = (\overline{\partial}_b h | g | \psi|^{-2})_{\widetilde{m},\hat{\theta}} = (\overline{\partial}_b (\psi^{-1} h) | \psi^{-1} g)_{\widetilde{m},\hat{\theta}}$$
$$= (\psi^{-1} h | \overline{\widetilde{\partial}}_b^{*,f} (\psi^{-1} g))_{\widetilde{m}} = (h | \psi \overline{\widetilde{\partial}}_b^{*,f} (\psi^{-1} g))_{m_1}.$$

Hence (5.5) follows. \Box

The next lemma clarifies the relation between $\overline{\partial}_{b,1}$ and $\widetilde{\overline{\partial}}_b$:

Lemma 5.2. We have

$$u \in \text{Dom}\,\overline{\partial}_{b,1} \text{ if and only if } \psi^{-1}u \in \text{Dom}\,\overline{\partial}_b.$$
 (5.6)

Moreover,

$$\overline{\partial}_{b,1}u = \psi \widetilde{\overline{\partial}}_b(\psi^{-1}u), \quad \forall u \in \text{Dom}\,\overline{\partial}_{b,1}.$$
(5.7)

Proof. Let $u \in \text{Dom }\overline{\partial}_{b,1}$. Then, there is an $h \in L^2_{(0,1)}(m_1, \hat{\theta})$ such that

$$(h \mid \alpha)_{m_1,\hat{\theta}} = (u \mid \overline{\partial}_{b,1}^{*,f} \alpha)_{m_1}, \quad \forall \alpha \in \Omega_0^{0,1}(X).$$

$$(5.8)$$

Note that $\overline{\partial}_{b,1}u := h$. From (5.5) and (5.8), it is easy to see that

$$(\psi^{-1}h \mid g)_{\widetilde{m},\widehat{\theta}} = (h \mid \psi g)_{m_1,\widehat{\theta}} = (u \mid \overline{\partial}_{b,1}^{*,f}(\psi g))_{m_1}$$
$$= (\psi^{-1}u \mid \psi^{-1}\overline{\partial}_{b,1}^{*,f}(\psi g))_{\widetilde{m}}$$
$$= (\psi^{-1}u \mid \widetilde{\overline{\partial}}_b^{*,f}g)_{\widetilde{m}}, \quad \forall g \in \Omega_0^{0,1}(X).$$
(5.9)

Since $\psi^{-1}h$ in $L^2_{(0,1)}(\widetilde{m}, \hat{\theta})$, from (5.9), we conclude that $\psi^{-1}u \in \text{Dom}\,\widetilde{\overline{\partial}}_b$ and

$$\widetilde{\overline{\partial}}_b(\psi^{-1}u) = \psi^{-1}h = \psi^{-1}\overline{\partial}_{b,1}u.$$

We have proved that if $u \in \text{Dom }\overline{\partial}_{b,1}$ then $\psi^{-1}u \in \text{Dom }\overline{\partial}_b$ and

$$\psi \overline{\overline{\partial}}_b(\psi^{-1}u) = \overline{\partial}_{b,1}u.$$

We can repeat the procedure above and conclude that if $v \in \text{Dom}\,\widetilde{\overline{\partial}}_b$ then $\psi v \in \text{Dom}\,\overline{\partial}_{b,1}$ and $\psi^{-1}\overline{\partial}_{b,1}(\psi v) = \widetilde{\overline{\partial}}_b v$. In fact, suppose $v \in \text{Dom}\,\widetilde{\overline{\partial}}_b$. It suffices to show that if $g \in \Omega_0^{0,1}(X)$, then

$$(\psi v \,|\, \overline{\partial}_{b,1}^{*,f} g)_{m_1} = (\psi \widetilde{\overline{\partial}}_b v \,|\, g)_{m_1,\hat{\theta}},$$

and that $\psi \tilde{\partial}_b v \in L^2(m_1, \hat{\theta})$. These are easy verifications, which we leave to the reader. The lemma follows. \Box

We have a corresponding lemma about the relation between $\overline{\partial}_{b,1}^*$ and $\overline{\overline{\partial}}_b^*$:

Lemma 5.3. We have

$$u \in \text{Dom}\,\overline{\partial}_{b,1}^* \text{ if and only if } \psi^{-1}u \in \text{Dom}\,\overline{\partial}_b.$$
 (5.10)

Moreover,

$$\overline{\partial}_{b,1}^* u = \psi \widetilde{\overline{\partial}}_b^* (\psi^{-1} u), \quad \forall u \in \operatorname{Dom} \overline{\partial}_{b,1}^*.$$
(5.11)

Proof. Let $u \in \text{Dom}\,\overline{\partial}_{b,1}^*$. We claim that $\psi^{-1}u \in \text{Dom}\,\overline{\partial}_b^*$ and $\overline{\partial}_b^*(\psi^{-1}u) = \psi^{-1}\overline{\partial}_{b,1}^*u$. Put $\overline{\partial}_{b,1}^*u = h \in L^2(m_1)$. By definition, we have

$$(\overline{\partial}_{b,1}g \mid u)_{m_1,\hat{\theta}} = (g \mid h)_{m_1}, \quad \forall g \in \text{Dom}\,\overline{\partial}_{b,1}.$$
(5.12)

From (5.2), Lemma 5.2 and (5.12), we have

$$(\overline{\partial}_{b}f \mid \psi^{-1}u)_{\widetilde{m},\widehat{\theta}} = (\psi\overline{\partial}_{b}f \mid u)_{m_{1},\widehat{\theta}} = (\overline{\partial}_{b,1}(\psi f) \mid u)_{m_{1},\widehat{\theta}}$$
$$= (\psi f \mid h)_{m_{1}} = (f \mid \psi^{-1}h)_{\widetilde{m}}, \quad \forall f \in \text{Dom}\,\widetilde{\partial}_{b}.$$
(5.13)

Thus, $\psi^{-1}u \in \text{Dom}\,\widetilde{\overline{\partial}}_b^*$ and $\widetilde{\overline{\partial}}_b^*(\psi^{-1}u) = \psi^{-1}h = \psi^{-1}\overline{\partial}_{b,1}^*u.$

We can repeat the procedure above and conclude that if $v \in \text{Dom}\,\widetilde{\overline{\partial}}_b^*$ then $\psi v \in \text{Dom}\,\overline{\partial}_{b,1}^*$ and $\overline{\partial}_{b,1}^*(\psi v) = \psi \widetilde{\overline{\partial}}_b^* v$. In fact, suppose $v \in \text{Dom}\,\widetilde{\overline{\partial}}_b^*$. It suffices to show that if $f \in \text{Dom}\,\overline{\partial}_{b,1}$, then

$$(\psi v \,|\, \overline{\partial}_{b,1} f)_{m_1,\hat{\theta}} = (\psi \widetilde{\overline{\partial}}_b^* v \,|\, f)_{m_1},$$

and $\psi \overline{\partial}_b^* v \in L^2(m_1)$. The verifications are easy and left to the reader. The lemma follows. \Box

Combining the above, we obtain (1.6) and (1.7):

Theorem 5.4. We have

$$u \in \text{Dom} \square_{b,1}$$
 if and only if $\frac{u}{\psi} \in \text{Dom} \widetilde{\square}_b$

and

$$\Box_{b,1} u = \psi \widetilde{\Box}_b(\frac{u}{\psi}), \quad \forall u \in \text{Dom} \, \Box_{b,1}.$$

Thus $\Box_{b,1}$ will have closed range in $L^2(m_1)$, if and only if $\tilde{\Box}_b$ has closed range in $L^2(\tilde{m})$. We will prove the latter in the next two sections, and that will establish Theorem 1.3.

6. The relation between $\widetilde{\Box}_b$ and $\hat{\Box}_b$

In the last section, we saw how the solution of $\Box_{b,1}$ is reduced to the solution of the intermediate operator $\widetilde{\Box}_b$. In this section, we see how the latter could be further reduced to solving \Box_b , which we introduced in Section 1.4. In particular, we will prove (1.8) and (1.9) there.

First, note that

$$L^{2}(\widetilde{m}) = L^{2}(\widehat{m}), \quad L^{2}_{(0,1)}(\widetilde{m},\widehat{\theta}) = L^{2}_{(0,1)}(\widehat{m},\widehat{\theta}).$$

In fact, from the expansions (1.5) of G_p and (4.14) of ψ , we see that near p,

$$\frac{\hat{m}}{\tilde{m}} = G_p^{-2} |\psi|^{-2} = 1 + a(z,t), \quad a(z,t) \in \mathcal{E}(\hat{\rho}^2).$$
(6.1)

Let

$$\hat{\overline{\partial}}_b^{*,f}: \Omega^{0,1}(X) \to C^\infty(X)$$

be the formal adjoint of $\overline{\partial}_b$ with respect to $(\cdot | \cdot)_{\hat{m}}, (\cdot | \cdot)_{\hat{m},\hat{\theta}}$. Then

$$\widetilde{\overline{\partial}}_{b}^{*,f} = \frac{\hat{m}}{\widetilde{m}} \hat{\overline{\partial}}_{b}^{*,f} \frac{\widetilde{m}}{\hat{m}},$$

so from the above, we see that

$$\widetilde{\overline{\partial}}_{b}^{*,f} = \widehat{\overline{\partial}}_{b}^{*,f} + g, \quad \text{for some } g \in \mathcal{E}(\hat{\rho}, T^{0,1}\hat{X}).$$
(6.2)

(Here we think of the (0,1) vector g as an element on the dual space of $\Lambda^{0,1}T^*\hat{X}$.) We then have the following lemma about the relation between $\tilde{\partial}_b$ and $\hat{\partial}_b$:

Lemma 6.1. We have $\tilde{\overline{\partial}}_b = \hat{\overline{\partial}}_b$. That is,

$$\operatorname{Dom} \widetilde{\overline{\partial}}_b = \operatorname{Dom} \hat{\overline{\partial}}_b, \quad and \quad \widetilde{\overline{\partial}}_b u = \hat{\overline{\partial}}_b u$$

for all u in the common domain of definition.

Proof. Let $u \in \text{Dom}\,\tilde{\overline{\partial}}_b$. We claim that $u \in \text{Dom}\,\hat{\overline{\partial}}_b$ and $\tilde{\overline{\partial}}_b u = \hat{\overline{\partial}}_b u$. By definition, there is an $h \in L^2_{(0,1)}(\tilde{m},\hat{\theta})$ such that

$$(h \mid \alpha)_{\widetilde{m},\widehat{\theta}} = (u \mid \widetilde{\overline{\partial}}_b^{*,f} \alpha)_{\widetilde{m}}, \quad \forall \alpha \in \Omega_0^{0,1}(X).$$

$$(6.3)$$

Note that $h = \tilde{\overline{\partial}}_b u$. From (6.2) and (6.3), we have

$$(h \mid \gamma)_{\hat{m},\hat{\theta}} = (h \mid \frac{\hat{m}}{\widetilde{m}}\gamma)_{\widetilde{m},\hat{\theta}} = (u \mid \overline{\partial}_{b}^{*,f}(\frac{\hat{m}}{\widetilde{m}}\gamma))_{\widetilde{m},\hat{\theta}} = (u \mid \frac{\hat{m}}{\widetilde{m}}\overline{\partial}_{b}^{*,f}\gamma)_{\widetilde{m},\hat{\theta}}$$
$$= (u \mid \overline{\partial}_{b}^{*,f}\gamma)_{\hat{m},\hat{\theta}}, \quad \forall \gamma \in \Omega_{0}^{0,1}(X).$$
(6.4)

Take $\alpha \in \Omega^{0,1}(\hat{X})$. Let

$$D_r := \left\{ x = (x_1, x_2, x_3) \in \mathbb{R}^3; |x|^2 := |x_1|^2 + |x_2|^2 + |x_3|^2 < r \right\}, \quad r > 0,$$

be a small ball. We identify D_r with a neighborhood of p. Take $\chi \in C_0^{\infty}(D_r, \overline{\mathbb{R}}_+), \chi = 1$ on $D_{\frac{r}{2}} := \left\{ x \in \mathbb{R}^3; |x|^2 < \frac{r}{2} \right\}$. Take $\epsilon > 0, \epsilon$ small and put $\alpha_{\epsilon} = (1 - \chi(\frac{x}{\epsilon}))\alpha \in \Omega_0^{0,1}(X)$. Then,

$$\begin{split} \|\alpha_{\epsilon} - \alpha\|_{\hat{m},\hat{\theta}}^{2} &= \int |\alpha|_{\hat{\theta}}^{2} \left|\chi(\frac{x}{\epsilon})\right|^{2} \hat{m} \leq C \int_{|x| \leq \epsilon r} \hat{m} \to 0 \text{ as } \epsilon \to 0, \\ \left\|\hat{\overline{\partial}}_{b}^{*,f}(\alpha_{\epsilon} - \alpha)\right\|_{\hat{m}}^{2} &\leq \int \left|\hat{\overline{\partial}}_{b}^{*,f}\alpha\right|_{\hat{m}}^{2} \left|\chi(\frac{x}{\epsilon})\right|^{2} \hat{m} + \int |\alpha|_{\hat{\theta}}^{2} \left|\frac{1}{\epsilon}(Z_{1}\chi)(\frac{x}{\epsilon})\right|^{2} \hat{m} \\ &\leq \frac{C_{1}}{\epsilon^{2}} \int_{|x| \leq \epsilon r} \hat{m} \to 0 \text{ as } \epsilon \to 0, \end{split}$$

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where C > 0, $C_1 > 0$ are constants independent of ϵ . We conclude that there exist $\alpha_j \in \Omega_0^{0,1}(X), j = 1, 2, \ldots$, such that

$$\lim_{j \to \infty} \|\alpha_j - \alpha\|_{\hat{m}, \hat{\theta}} = 0, \quad \lim_{j \to \infty} \left\| \hat{\bar{\partial}}_b^{*, f}(\alpha_j - \alpha) \right\|_{\hat{m}} = 0.$$

Combining this with (6.4), we have

$$(h \mid \alpha)_{\hat{m},\hat{\theta}} = \lim_{j \to \infty} (h \mid \alpha_j)_{\hat{m},\hat{\theta}} = \lim_{j \to \infty} (u \mid \hat{\overline{\partial}}_b^{*,f} \alpha_j)_{\hat{m}} = (u \mid \hat{\overline{\partial}}_b^{*,f} \alpha)_{\hat{m}}.$$

Thus, $u \in \text{Dom}\,\hat{\overline{\partial}}_b$ and $\hat{\overline{\partial}}_b u = \widetilde{\overline{\partial}}_b u = h$.

We have proved that if $u \in \text{Dom}\,\tilde{\overline{\partial}}_b$ then $u \in \text{Dom}\,\hat{\overline{\partial}}_b$ and $\hat{\overline{\partial}}_b u = \widetilde{\overline{\partial}}_b u$.

We can repeat the procedure above and conclude that if $v \in \text{Dom}\,\overline{\partial}_b$ then $v \in \text{Dom}\,\overline{\partial}_b$ and $\widetilde{\partial}_b v = \overline{\partial}_b v$. In fact, suppose $v \in \text{Dom}\,\overline{\partial}_b$. It suffices to show that if $\gamma \in \Omega_0^{0,1}(X)$, then

$$(v \,|\, \widehat{\overline{\partial}}_b^{*,f} \gamma)_{\hat{m}} = (\widetilde{\overline{\partial}}_b v \,|\, \gamma)_{\hat{m},\hat{\theta}},$$

which is easy to verify. The lemma follows. \Box

Next, we want to understand the relation between $\tilde{\overline{\partial}}_b^*$ and $\hat{\overline{\partial}}_b^*$. To do so, we need the following lemma:

Lemma 6.2. Let $v \in \text{Dom}\,\hat{\bar{\partial}}_b$. Then, $\underline{\widetilde{m}}_{\hat{m}}v \in \text{Dom}\,\hat{\bar{\partial}}_b$ and

$$\hat{\overline{\partial}}_b(\frac{\widetilde{m}}{\hat{m}}v) = \frac{\widetilde{m}}{\hat{m}}\hat{\overline{\partial}}_bv - \frac{\widetilde{m}}{\hat{m}}g^*v$$

where g^* is the (0,1) form dual to the (0,1) vector g with respect to the Hermitian inner product $\langle \cdot | \cdot \rangle_{\hat{\theta}}$ on $T^{0,1}X$.

Proof. For all $\alpha \in \Omega_0^{0,1}(X)$, we have

$$\begin{split} (\frac{\widetilde{m}}{\widehat{m}}v \mid \hat{\overline{\partial}}_{b}^{*,f}\alpha)_{\widehat{m}} &= (\frac{\widetilde{m}}{\widehat{m}}v \mid (\widetilde{\overline{\partial}}_{b}^{*,f} - g)\alpha)_{\widehat{m}} \text{ here we used (6.2)} \\ &= (v \mid (\widetilde{\overline{\partial}}_{b}^{*,f} - g)\alpha)_{\widetilde{m}} \\ &= (\hat{\overline{\partial}}_{b}v \mid \alpha)_{\widetilde{m},\widehat{\theta}} - (g^{*}v \mid \alpha)_{\widetilde{m},\widehat{\theta}} \text{ here we used Lemma 6.1} \\ &= (\frac{\widetilde{m}}{\widehat{m}}\hat{\overline{\partial}}_{b}v - \frac{\widetilde{m}}{\widehat{m}}g^{*}v \mid \alpha)_{\widehat{m},\widehat{\theta}}. \end{split}$$

We can now repeat the procedure in the proof of Lemma 6.1 and conclude that

$$(\frac{\widetilde{m}}{\widehat{m}}v\mid\hat{\overline{\partial}}_{b}^{*,f}h)_{\widehat{m}}=(\frac{\widetilde{m}}{\widehat{m}}\hat{\overline{\partial}}_{b}v-\frac{\widetilde{m}}{\widehat{m}}g^{*}v\mid h)_{\widehat{m},\widehat{\theta}}, \quad \forall h\in\Omega^{0,1}(\hat{X}).$$

The lemma follows. \Box

We can now prove:

Lemma 6.3. We have $\operatorname{Dom} \hat{\overline{\partial}}_b^* = \operatorname{Dom} \tilde{\overline{\partial}}_b^*$ and

$$\tilde{\bar{\partial}}_{b}^{*}u = \hat{\bar{\partial}}_{b}^{*}u + gu, \quad \forall u \in \text{Dom}\,\hat{\bar{\partial}}_{b}^{*} = \text{Dom}\,\tilde{\bar{\partial}}_{b}^{*}.$$
(6.5)

Proof. Let $u \in \text{Dom}\,\hat{\overline{\partial}}_b^*$. We claim that $u \in \text{Dom}\,\widetilde{\overline{\partial}}_b^*$ and $\widetilde{\overline{\partial}}_b^* u = \hat{\overline{\partial}}_b^* u + gu$. From Lemma 6.1 and Lemma 6.2, we can check that for every $v \in \text{Dom}\,\widetilde{\overline{\partial}}_b$,

$$\begin{aligned} (v \mid (\hat{\overline{\partial}}_{b}^{*} + g)u)_{\widetilde{m}} &= (\frac{\widetilde{m}}{\widehat{m}}v \mid \hat{\overline{\partial}}_{b}^{*}u)_{\widehat{m}} + (v \mid gu)_{\widetilde{m}} \\ &= (\hat{\overline{\partial}}_{b}(\frac{\widetilde{m}}{\widehat{m}}v) \mid u)_{\widehat{m},\widehat{\theta}} + (g^{*}v \mid u)_{\widetilde{m},\widehat{\theta}} \\ &= (\frac{\widetilde{m}}{\widehat{m}}\hat{\overline{\partial}}_{b}v - \frac{\widetilde{m}}{\widehat{m}}g^{*}v \mid u)_{\widehat{m},\widehat{\theta}} + (\frac{\widetilde{m}}{\widehat{m}}g^{*}v \mid u)_{\widehat{m}} \\ &= (\hat{\overline{\partial}}_{b}v \mid u)_{\widetilde{m},\widehat{\theta}} = (\tilde{\overline{\partial}}_{b}v \mid u)_{\widetilde{m},\widehat{\theta}}. \end{aligned}$$

Thus, $u \in \text{Dom}\,\widetilde{\overline{\partial}}_b^*$ and $\widetilde{\overline{\partial}}_b^* u = \hat{\overline{\partial}}_b^* u + gu$.

Similarly, for $u \in \text{Dom}\,\widetilde{\overline{\partial}}_b^*$, we can repeat the procedure above and conclude that $u \in \text{Dom}\,\widehat{\overline{\partial}}_b^*$. In fact, suppose $u \in \text{Dom}\,\widetilde{\overline{\partial}}_b^*$. It suffices to show that for every $v \in \text{Dom}\,\widehat{\overline{\partial}}_b$,

$$(u\,|\,\widehat{\overline{\partial}}_b v)_{\hat{m},\hat{\theta}} = (\widetilde{\overline{\partial}}_b^* u - gu\,|\,v)_{\hat{m}},$$

and that $\tilde{\overline{\partial}}_b^* u - gu \in L^2(\hat{m})$. This can be proved in exactly the same manner as above. The lemma follows. \Box

It follows from the above that (1.8) and (1.9) hold:

Theorem 6.4.

$$\operatorname{Dom} \tilde{\Box}_b = \operatorname{Dom} \hat{\Box}_b,$$

and

$$\widetilde{\Box}_b u = \hat{\Box}_b u + g \overline{\partial}_b u, \quad \forall u \in \operatorname{Dom} \widetilde{\Box}_b.$$

This allows one to understand solutions to $\widehat{\Box}_b$ via solutions to $\widehat{\Box}_b$, as we will see below.

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7. Proof of Theorem 1.3

In this section, we will see that \square_b and $\square_{b,1}$ have closed ranges in $L^2(\widetilde{m})$ and $L^2(m_1)$ respectively. The analogous property for \square_b is well-known by the CR embeddability of \widehat{X} into \mathbb{C}^N .

Theorem 7.1. The operator

$$\widetilde{\overline{\partial}}_b : \operatorname{Dom} \widetilde{\overline{\partial}}_b \subset L^2(\widetilde{m}) \to L^2_{(0,1)}(\widetilde{m}, \widehat{\theta})$$

has closed range.

Proof. One simply notes that since \hat{X} is CR embeddable in some \mathbb{C}^N , by the result of [17], $\hat{\bar{\partial}}_b$: Dom $\hat{\bar{\partial}}_b \subset L^2(\hat{m}) \to L^2_{(0,1)}(\hat{m}, \hat{\theta})$ has closed range in $L^2_{(0,1)}(\hat{m}, \hat{\theta})$. By the identity of $\hat{\bar{\partial}}_b$ with $\tilde{\bar{\partial}}_b$ as in Lemma 6.1, it follows that $\tilde{\bar{\partial}}_b$ has closed range in $L^2_{(0,1)}(\hat{m}, \hat{\theta})$. \Box

It is now a standard matter to prove:

Theorem 7.2. The operator

$$\widetilde{\Box}_b$$
: Dom $\widetilde{\Box}_b \subset L^2(\widetilde{m}) \to L^2(\widetilde{m})$

has closed range.

Proof. By Theorem 7.1, there is a constant c > 0 such that

$$\left\|\overline{\widetilde{\partial}}_{b}u\right\|_{\widetilde{m},\widehat{\theta}}^{2} \ge c \left\|u\right\|_{\widetilde{m}}^{2}, \quad \forall u \bot \operatorname{Ker}\widetilde{\widetilde{\partial}}_{b}.$$

$$(7.1)$$

Let $f \in \text{Dom}\,\widetilde{\Box}_b \bigcap (\text{Ker}\,\widetilde{\Box}_b)^{\perp}$. It is not difficult to see that $\text{Ker}\,\widetilde{\Box}_b = \text{Ker}\,\widetilde{\overline{\partial}}_b$. Thus, $f \in \text{Dom}\,\widetilde{\overline{\partial}}_b \bigcap (\text{Ker}\,\widetilde{\overline{\partial}}_b)^{\perp}$. From this observation and (7.1), we have

$$\|\widetilde{\Box}_b f\|_{\widetilde{m}} \|f\|_{\widetilde{m}} \ge (\widetilde{\Box}_b f \mid f)_{\widetilde{m}} = (\widetilde{\overline{\partial}}_b f \mid \widetilde{\overline{\partial}}_b f)_{\widetilde{m},\hat{\theta}} \ge c \|f\|_{\widetilde{m}}^2,$$

where c > 0 is the constant as in (7.1). Thus,

$$\|\widetilde{\Box}_b f\|_{\widetilde{m}} \ge c \, \|f\|_{\widetilde{m}} \,, \quad \forall f \in \operatorname{Dom} \widetilde{\Box}_b \bigcap (\operatorname{Ker} \widetilde{\Box}_b)^{\perp}, \tag{7.2}$$

where c > 0 is the constant as in (7.1). From (7.2), it is easy to see that $\widetilde{\Box}_b$ has closed range. The theorem follows. \Box

Now we use the closed range property of $\tilde{\overline{\partial}}_b$ to prove the same for $\bar{\partial}_{b,1}$.

Theorem 7.3. The operator

$$\overline{\partial}_{b,1} : \operatorname{Dom} \overline{\partial}_{b,1} \subset L^2(m_1) \to L^2_{(0,1)}(m_1, \hat{\theta})$$

has closed range.

Proof. Let $f_j \in \text{Dom }\overline{\partial}_{b,1}, j = 1, 2, ..., \overline{\partial}_{b,1}f_j = g_j \in L^2_{(0,1)}(m_1, \hat{\theta}), j = 1, 2, ...$ We assume that there is a function $g \in L^2_{(0,1)}(m_1, \hat{\theta})$ such that $\lim_{j\to\infty} \|g_j - g\|_{m_1,\hat{\theta}} = 0$. We are going to show that $g \in \text{Ran }\overline{\partial}_{b,1}$. From (5.6), (5.7) and (5.2), we see that

$$\psi^{-1} f_{j} \in \text{Dom}\,\tilde{\bar{\partial}}_{b}, \quad j = 1, 2, \dots,$$

$$\tilde{\bar{\partial}}_{b}(\psi^{-1} f_{j}) = \psi^{-1} g_{j} \in L^{2}_{(0,1)}(\tilde{m}, \hat{\theta}), \quad j = 1, 2, \dots,$$

$$\lim_{j \to \infty} \left\| \psi^{-1} g_{j} - \psi^{-1} g \right\|_{\tilde{m}, \hat{\theta}} = 0.$$
 (7.3)

Since $\overline{\partial}_b$ has closed range, we can find $\widetilde{h} \in \text{Dom}\,\overline{\partial}_b$ such that $\overline{\partial}_b \widetilde{h} = \psi^{-1}g$. Put $h = \psi \widetilde{h} \in L^2(m_1)$. From (5.6) and (5.7), we see that $h \in \text{Dom}\,\overline{\partial}_{b,1}$ and $\overline{\partial}_{b,1}h = g$. Thus, $g \in \text{Ran}\,\overline{\partial}_{b,1}$. The theorem follows. \Box

From Theorem 7.3, we can repeat the proof of Theorem 7.2 and conclude the proof of Theorem 1.3.

8. Proof of Theorem 1.4

In the last section, we have seen that $\Box_{b,1}$ and $\widetilde{\Box}_b$ have closed ranges in L^2 . Thus one can define the partial inverses N and \widetilde{N} of $\Box_{b,1}$ and $\widetilde{\Box}_b$ respectively (cf. Section 1.5). Furthermore, we write Π and $\widetilde{\Pi}$ for the Szegö projections, which are orthogonal projections onto Ker $\widetilde{\Box}_b$ and Ker $\Box_{b,1}$ respectively, as in Sections 1.3 and 1.4. Our goal is to understand N and Π , as in Theorem 1.4. But from (5.1), (5.2) and Theorem 5.4, we obtain

$$\Box_{b,1}\psi \widetilde{N}\frac{1}{\psi} + \psi \widetilde{\Pi}\frac{1}{\psi} = I \quad \text{on } L^2(m_1).$$

Thus we obtain (1.10), namely

$$N = \psi \widetilde{N} \frac{1}{\psi}$$
 and $\Pi = \psi \widetilde{\Pi} \frac{1}{\psi}$.

The analysis of N and Π then reduces to the analysis of \widetilde{N} and $\widetilde{\Pi}$; in fact, to prove Theorem 1.4, it suffices to prove instead (1.11) and (1.12):

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Theorem 8.1. $\widetilde{\Pi}$ and \widetilde{N} extend as continuous operators

$$\widetilde{\Pi}: \mathcal{E}(\hat{\rho}^{-4+\delta}) \to \mathcal{E}(\hat{\rho}^{-4+\delta})$$
$$\widetilde{N}: \mathcal{E}(\hat{\rho}^{-4+\delta}) \to \mathcal{E}(\hat{\rho}^{-2+\delta})$$

for every $0 < \delta < 2$.

We will achieve this by reducing to the analogous properties of \hat{N} and $\hat{\Pi}$, which we proved in Section 2.

The starting point is the following lemma:

Lemma 8.2. On $L^2(\widetilde{m})$, we have

$$\widetilde{\Pi}(I+\hat{R}) = \widehat{\Pi},\tag{8.1}$$

$$\widetilde{N}(I+\hat{R}) = (I-\widetilde{\Pi})\hat{N},\tag{8.2}$$

where $\hat{R} = g \hat{\overline{\partial}}_b \hat{N}$.

Proof. First, we know that $\hat{N} : L^2(\hat{m}) \to \text{Dom }\hat{\Box}_b = \text{Dom }\widetilde{\Box}_b$. From Theorem 6.4, we can check that

$$\widetilde{\Box}_b \hat{N} + \hat{\Pi} = \hat{\Box}_b \hat{N} + \hat{\Pi} + \hat{R} = I + \hat{R}.$$

From this, we have

$$\widetilde{\Pi}(I+\hat{R}) = \widetilde{\Pi}(\widetilde{\Box}_b \hat{N} + \hat{\Pi}) = \widetilde{\Pi}\hat{\Pi}.$$
(8.3)

On the other hand, we have

$$\hat{\Pi} = (\widetilde{N}\widetilde{\Box}_b + \widetilde{\Pi})\hat{\Pi} = \widetilde{\Pi}\hat{\Pi}.$$
(8.4)

From (8.3) and (8.4), we get (8.1).

Now, from Theorem 6.4 again, we have

$$\hat{N} = (\hat{N} \square_b + \Pi) \hat{N}$$

= $(\tilde{N} \square_b + \Pi) \hat{N} + \tilde{N} \hat{R}$
= $\tilde{N} (I - \Pi) + \Pi \hat{N} + \tilde{N} \hat{R}$
= $\tilde{N} + \Pi \hat{N} + \tilde{N} \hat{R};$

in the last line we used that fact that $\widetilde{N}\widehat{\Pi} = 0$. (8.2) then follows, and we are done. \Box

We now extend the definitions of Π and \tilde{N} to $\mathcal{E}(\hat{\rho}^{-4+\delta})$, $0 < \delta < 2$. The problem is that one does not have a good inverse for $(I + \hat{R})$ in (8.1) and (8.2). The key then is to rewrite (8.1) and (8.2) as

$$\tilde{\Pi} = \hat{\Pi} - \tilde{\Pi}\hat{R},\tag{8.5}$$

$$\widetilde{N} = (I - \widetilde{\Pi})\hat{N} - \widetilde{N}\hat{R}.$$
(8.6)

Note that \hat{R} extends to a continuous operator

$$\hat{R}: \mathcal{E}(\hat{\rho}^{-4+\delta}) \to \mathcal{E}(\hat{\rho}^{-2+\delta}) \subset L^2(\widetilde{m})$$

for every $0 < \delta < 2$, since $\hat{R} = g\hat{\bar{\partial}}_b \hat{N}$, and \hat{N} satisfies the analogous property. Thus $\Pi \hat{R}$ maps $\mathcal{E}(\hat{\rho}^{-4+\delta})$ continuously to $L^2(\tilde{m})$. Since Π maps $\mathcal{E}(\hat{\rho}^{-4+\delta})$ continuously to $\mathcal{E}(\hat{\rho}^{-4+\delta})$, by (8.5), we have extended the domain of definition of Π to $\mathcal{E}(\hat{\rho}^{-4+\delta})$. Furthermore, $\tilde{N}\hat{R}$ maps $\mathcal{E}(\hat{\rho}^{-4+\delta})$ continuously to $L^2(\tilde{m})$, and \hat{N} maps $\mathcal{E}(\hat{\rho}^{-4+\delta})$ continuously to $\mathcal{E}(\hat{\rho}^{-2+\delta}) \subset L^2(\tilde{m})$. Thus together with the continuity of Π on $L^2(\tilde{m})$, from (8.6), we see that \tilde{N} extends as a continuous map $\mathcal{E}(\hat{\rho}^{-4+\delta}) \to L^2(\tilde{m})$.

To proceed further, let's write $\hat{\Pi}^{*,\widetilde{m}}$, $\hat{N}^{*,\widetilde{m}}$ and $\hat{R}^{*,\widetilde{m}}$ for the adjoints of $\hat{\Pi}$, \hat{N} and \hat{R} with respect to the inner product of $L^2(\widetilde{m})$. We note that

$$\begin{split} \hat{\Pi}^{*,\widetilde{m}} &= \frac{\hat{m}}{\widetilde{m}}\hat{\Pi}\frac{\widetilde{m}}{\hat{m}}, \\ \hat{N}^{*,\widetilde{m}} &= \frac{\hat{m}}{\widetilde{m}}\hat{N}\frac{\widetilde{m}}{\hat{m}}, \\ \hat{R}^{*,\widetilde{m}} &= \frac{\hat{m}}{\widetilde{m}}\hat{N}\hat{\overline{\partial}}_{b}^{*}(g^{*}\frac{\widetilde{m}}{\hat{m}}). \end{split}$$

where $\tilde{m}/\hat{m} := G_p^2 |\psi|^2$ is the density of \tilde{m} with respect to \hat{m} , and similarly $\hat{m}/\tilde{m} := G_p^{-2} |\psi|^{-2}$. Here g^* is the (0,1) form dual to g. Since $\tilde{m}/\hat{m}, \hat{m}/\tilde{m} \in \mathcal{E}(\hat{\rho}^0)$, one can show that

$$\hat{\Pi}^{*,\widetilde{m}}: \mathcal{E}(\hat{\rho}^{-4+\delta}) \to \mathcal{E}(\hat{\rho}^{-4+\delta}) \tag{8.7}$$

$$\hat{N}^{*,\widetilde{m}}: \mathcal{E}(\hat{\rho}^{-4+\delta}) \to \mathcal{E}(\hat{\rho}^{-2+\delta})$$
(8.8)

$$\hat{R}^{*,\widetilde{m}}: \mathcal{E}(\hat{\rho}^{-4+\delta}) \to \mathcal{E}(\hat{\rho}^{-2+\delta})$$
(8.9)

for every $0 < \delta < 2$; these are easy consequences of the analogous properties of $\hat{\Pi}$, \hat{N} and \hat{R} . The problem is that it is not clear that $(I + \hat{R}^{*,\widetilde{m}})$ is invertible on $\mathcal{E}(\hat{\rho}^{-4+\delta})$; if it is, then we can invoke (1.18) and (1.19) and our proof of Theorem 1.4 would be much easier. In order to get around this problem, we introduce a cut-off χ , as was explained in Section 1.4; we will prove

Theorem 8.3. Let $\chi \in C^{\infty}(\hat{X})$ with $\chi = 1$ near p. Then

$$(1-\chi)\widetilde{N}, (1-\chi)\widetilde{\Pi}: \mathcal{E}(\hat{\rho}^{-4+\delta}) \to C_0^{\infty}(X)$$

are continuous for $0 < \delta < 2$.

Theorem 8.4. If the support of $\chi \in C^{\infty}(\hat{X})$ is a sufficiently small neighborhood of p, then $(I + \hat{R}^{*,\widetilde{m}}\chi)$ is invertible on $L^{2}(\widetilde{m})$, and extends to a linear map

$$(I + \hat{R}^{*,\widetilde{m}}\chi)^{-1}: \mathcal{E}(\hat{\rho}^{-4+\delta}) \to \mathcal{E}(\hat{\rho}^{-4+\delta})$$

for every $0 < \delta < 4$.

Assuming these for the moment. Then one can finish the proof of Theorem 8.1 using (1.23) and (1.24) as explained in Section 1.4 shortly after these identities. Theorem 1.4 then follows from (1.10) and Theorem 8.1. We omit the details.

We now turn to the proofs of the theorems above.

Proof of Theorem 8.4. First, observe that

$$I + \hat{R}^{*,\widetilde{m}}\chi = \frac{\hat{m}}{\widetilde{m}}(I + \hat{N}\hat{\bar{\partial}}_{b}^{*}\chi g^{*})\frac{\widetilde{m}}{\hat{m}};$$

and \tilde{m}/\hat{m} , $\hat{m}/\tilde{m} \in \mathcal{E}(\hat{\rho}^0)$. Thus it suffices to prove that $I + \hat{N}\hat{\bar{\partial}}_b^* \chi g^*$ is invertible on $L^2(\tilde{m})$, and extends to a linear map

$$(I + \hat{N}\hat{\bar{\partial}}_b^* \chi g^*)^{-1} : \mathcal{E}(\hat{\rho}^{-4+\delta}) \to \mathcal{E}(\hat{\rho}^{-4+\delta})$$

for every $0 < \delta < 4$. But this follows from Theorem 3.1, since $\chi g^* \in \mathcal{E}(\hat{\rho}, \Lambda^{0,1}T^*\hat{X})$ has compact support in a sufficiently small neighborhood of p. Thus we are done. \Box

Proof of Theorem 8.3. We need to recall the Kohn Laplacian $\hat{\Box}_{b,\varepsilon} := \hat{\bar{\partial}}_{b,\varepsilon}^* \hat{\bar{\partial}}_{b,\varepsilon}$ with respect to the volume form $m_{\varepsilon} := \eta_{\varepsilon} \hat{m} + (1 - \eta_{\varepsilon}) \tilde{m}$ as described in Section 1.4. First

$$\frac{\hat{m}_{\varepsilon}}{\tilde{m}} = \eta_{\varepsilon} G_p^{-2} |\psi|^{-2} + (1 - \eta_{\varepsilon}) = \eta_{\varepsilon} (1 + a) + (1 - \eta_{\varepsilon}) = 1 + \eta_{\varepsilon} a.$$

The upshot is that $\eta_{\varepsilon} a \in \mathcal{E}(\hat{\rho}^2, T^{0,1}\hat{X})$ has compact support near p. Thus if we follow the construction in Section 6, there will exist some $g_{\varepsilon} \in \mathcal{E}(\hat{\rho}^1, T^{0,1}\hat{X})$ (possibly non-smooth near p) such that

$$\widetilde{\Box}_{b}u = \hat{\Box}_{b,\varepsilon}u + g_{\varepsilon}\overline{\bar{\partial}}_{b}u, \quad \forall u \in \text{Dom}\,\tilde{\Box}_{b};$$
(8.10)

in addition g_{ε} will be compactly supported in the support of η_{ε} . It is known that $\hat{\Box}_{b,\varepsilon}$ has closed range in $L^2(\hat{m}_{\varepsilon})$. Thus one can define the partial inverse \hat{N}_{ε} of $\hat{\Box}_{b,\varepsilon}$, as well as the Szegö projection $\hat{\Pi}_{\varepsilon}$ onto the kernel of $\hat{\Box}_{b,\varepsilon}$, such that

$$\hat{\Box}_{b,\varepsilon}\hat{N}_{\varepsilon} + \hat{\Pi}_{\varepsilon} = I.$$

One can then repeat the proof of Lemma 8.2, and show that

$$\widetilde{\Pi}(I + \hat{R}_{\varepsilon}) = \hat{\Pi}_{\varepsilon} \tag{8.11}$$

$$\widetilde{N}(I + \hat{R}_{\varepsilon}) = (I - \widetilde{\Pi})\hat{N}_{\varepsilon}$$
(8.12)

on $L^2(\widetilde{m})$, where

 $\hat{R}_{\varepsilon} := g_{\varepsilon} \hat{\overline{\partial}}_{b,\varepsilon} \hat{N}_{\varepsilon}.$

Now we write $\hat{\Pi}_{\varepsilon}^{*,\widetilde{m}}$, $\hat{N}_{\varepsilon}^{*,\widetilde{m}}$ and $\hat{R}_{\varepsilon}^{*,\widetilde{m}}$ for the adjoints of $\hat{\Pi}_{\varepsilon}$, \hat{N}_{ε} and \hat{R}_{ε} with respect to $L^{2}(\widetilde{m})$. By taking adjoints of (8.11) and (8.12), and multiplying by $(1 - \chi)$, it follows that

$$(1-\chi)\widetilde{\Pi} = (1-\chi)\widehat{\Pi}_{\varepsilon}^{*,\widetilde{m}} - (1-\chi)\widehat{R}_{\varepsilon}^{*,\widetilde{m}}\widetilde{\Pi}, \qquad (8.13)$$

$$(1-\chi)\widetilde{N} = (1-\chi)\hat{N}_{\varepsilon}^{*,\widetilde{m}}(I-\widetilde{\Pi}) - (1-\chi)\hat{R}_{\varepsilon}^{*,\widetilde{m}}\widetilde{N}.$$
(8.14)

But

$$(1-\chi)\hat{R}^{*,\widetilde{m}}_{\varepsilon} = \frac{\hat{m}_{\varepsilon}}{\widetilde{m}}(1-\chi)\hat{N}_{\varepsilon}\hat{\overline{\partial}}^{*}_{b,\varepsilon}g^{*}_{\varepsilon}\frac{\widetilde{m}}{\hat{m}_{\varepsilon}},$$

and if ε is chosen sufficiently small (so that the support of g_{ε} is disjoint from that of $1-\chi$), then $(1-\chi)\hat{N}_{\varepsilon}\hat{\bar{\partial}}_{b,\varepsilon}^{*}g_{\varepsilon}^{*}$ is an infinitely smoothing pseudodifferential operator, by pseudolocality of \hat{N}_{ε} . Hence the last term of (8.13), and also the last term of (8.14), map $\mathcal{E}(\hat{\rho}^{-4+\delta})$ into $C_{0}^{\infty}(X)$. Since for every $0 < \varepsilon < 1$ and every $0 < \delta < 2$,

$$\hat{\Pi}_{\varepsilon}^{*,\widetilde{m}} = \frac{\hat{m}_{\varepsilon}}{\widetilde{m}} \hat{\Pi}_{\varepsilon} \frac{\widetilde{m}}{\hat{m}_{\varepsilon}} : \mathcal{E}(\hat{\rho}^{-4+\delta}) \to \mathcal{E}(\hat{\rho}^{-4+\delta}),$$

and

$$\hat{N}_{\varepsilon}^{*,\widetilde{m}} = \frac{\hat{m}_{\varepsilon}}{\widetilde{m}} \hat{N}_{\varepsilon} \frac{\widetilde{m}}{\hat{m}_{\varepsilon}} : \mathcal{E}(\hat{\rho}^{-4+\delta}) \to \mathcal{E}(\hat{\rho}^{-2+\delta}),$$

it follows from (8.13) that $(1-\chi)\widetilde{\Pi}$ maps $\mathcal{E}(\hat{\rho}^{-4+\delta})$ continuously into $C_0^{\infty}(X)$ as desired. This then implies the corresponding result for $(1-\chi)\widetilde{N}$ by (8.14), and we are done. \Box

9. Proof of Theorem 1.5

In this section, we will complete the proof of Theorem 1.5. To begin with, we have the following lemma:

Lemma 9.1. If $\alpha \in \mathcal{E}(\hat{\rho}^0, \Lambda^{0,1}T^*\hat{X})$, then $\overline{\partial}_{b,1}^* \alpha \in \mathcal{E}(\hat{\rho}^{-1})$, and

$$\Pi \overline{\partial}_{b,1}^* \alpha = 0.$$

Proof. Given $\alpha \in \mathcal{E}(\hat{\rho}^0, \Lambda^{0,1}T^*\hat{X})$, and any $\gamma \in (0, 1)$, there exists a sequence of smooth and compactly supported $\alpha_j \in \Omega_0^{0,1}(X)$ such that

$$\alpha_j \to \alpha \quad \text{in } \mathcal{E}(\hat{\rho}^{-\gamma}, \Lambda^{0,1}T^*\hat{X}).$$

Then

$$\overline{\partial}_{b,1}^* \alpha_j \to \overline{\partial}_{b,1}^* \alpha \quad \text{in } \mathcal{E}(\hat{\rho}^{-1-\gamma}, \Lambda^{0,1}T^*\hat{X}),$$

so by continuity of Π on $\mathcal{E}(\hat{\rho}^{-1-\gamma}, \Lambda^{0,1}T^*\hat{X})$, we have

$$\Pi \overline{\partial}_{b,1}^* \alpha = \lim_{j \to \infty} \Pi \overline{\partial}_{b,1}^* \alpha_j.$$

But the right hand side here is zero, since

$$\Pi \overline{\partial}_{b,1}^* = 0 \quad \text{on Dom } \overline{\partial}_{b,1}^*$$

and $\alpha_j \in \text{Dom }\overline{\partial}_{b,1}^*$ for all j. Hence we are done. \Box

Now to prove of Theorem 1.5, suppose $f = G_p^2 F$ and $F = \Box_b \tilde{\beta}$ as in Section 1.3, where

$$\widetilde{\beta} = \beta_0 + \beta_1, \quad \beta_0 = \chi(z,t) \frac{i\overline{z}}{|z|^2 - it} \in \mathcal{E}(\hat{\rho}^{-1}), \quad \beta_1 \in \mathcal{E}(\hat{\rho}^1).$$

Then $f = \Box_{b,1} \tilde{\beta}$. Our goal is to compute $\Pi f = \Pi \overline{\partial}_{b,1}^* (\overline{\partial}_{b,1} \tilde{\beta})$. The problem is that $\overline{\partial}_{b,1} \tilde{\beta}$ is not in $\mathcal{E}(\hat{\rho}^0, \Lambda^{0,1} T^* \hat{X})$; otherwise we could simply apply the above lemma to conclude. Nevertheless, we will write $\overline{\partial}_{b,1} \tilde{\beta}$ as the sum of a main term and an error, where the error is in $\mathcal{E}(\hat{\rho}^0, \Lambda^{0,1} T^* \hat{X})$, and the $\overline{\partial}_{b,1}^*$ of the main term can be approximated in $\mathcal{E}(\hat{\rho}^{-1})$ by the $\overline{\partial}_{b,1}^*$ of some forms in $\mathcal{E}(\hat{\rho}^0, \Lambda^{0,1} T^* \hat{X})$. Then we can conclude using the lemma, and the continuity of Π on $\mathcal{E}(\hat{\rho}^{-1})$.

To begin with, note that

$$\overline{\partial}_{b,1}\widetilde{\beta} = \chi \overline{\partial}_{b,1} \frac{i\overline{z}}{|z|^2 - it} + (\overline{\partial}_{b,1}\chi) \frac{i\overline{z}}{|z|^2 - it} + \overline{\partial}_{b,1}\beta_1,$$

and the last two terms are in $\mathcal{E}(\hat{\rho}^0, \Lambda^{0,1}T^*\hat{X})$. Thus the key is to compute the first term. Suppose we pick a local section \hat{Z}_1 of $T^{1,0}\hat{X}$ near p, with $\langle \hat{Z}_1|\hat{Z}_1\rangle_{\hat{\mu}} = 1$, such that \hat{Z}_1 admits the expansion (1.31) in CR normal coordinates (z,t) near p. We also write $\hat{Z}_{\bar{1}}$ for $\overline{\hat{Z}_1}$, and $\hat{Z}^{\bar{1}}$ for the dual (0,1) form of $\hat{Z}_{\bar{1}}$. Then by the expansion (1.31) of \hat{Z}_1 , we have

$$\chi \overline{\partial}_{b,1} \frac{i\overline{z}}{|z|^2 - it} = \chi \hat{Z}_{\bar{1}} \left(\frac{i\overline{z}}{|z|^2 - it} \right) \hat{Z}^{\bar{1}} = -i\chi \frac{|z|^2 + it}{(|z|^2 - it)^2} \hat{Z}^{\bar{1}} + \text{error},$$

where the error is in $\mathcal{E}(\hat{\rho}^0, \Lambda^{0,1}T^*\hat{X})$. Thus

$$\overline{\partial}_{b,1}\widetilde{\beta} = -i\chi \frac{|z|^2 + it}{(|z|^2 - it)^2} \hat{Z}^{\bar{1}} + \text{error},$$

where the first term is in $\mathcal{E}(\hat{\rho}^{-2}, \Lambda^{0,1}T^*\hat{X})$, and the error is in $\mathcal{E}(\hat{\rho}^0, \Lambda^{0,1}T^*\hat{X})$.

To proceed further, we write

$$m_1 = G_p^2 \hat{\theta} \wedge d\hat{\theta} = 2iv(z,t)dt \wedge d\bar{z} \wedge dz$$

for some function v(z,t) near p. Then by the expansion (1.5) of G_p and the expansion (1.33) of $\hat{\theta}$, we have

$$v(z,t) = \frac{1}{4\pi^2 \hat{\rho}^4} + \frac{A}{\pi \hat{\rho}^2} + \text{error in } \mathcal{E}(\hat{\rho}^{-1}).$$

Now we define a (0, 1) form

$$\alpha_0 := -2\pi i \chi v^{-1} \overline{\psi}^{-3} \hat{Z}^{\bar{1}}$$

near p. Then using the expansion (4.14) of ψ , we get

$$-i\chi \frac{|z|^2 + it}{(|z|^2 - it)^2} \hat{Z}^{\bar{1}} = \alpha_0 + \text{error in } \mathcal{E}(\hat{\rho}^0, \Lambda^{0,1} T^* \hat{X})$$

It follows that

$$\overline{\partial}_{b,1}\widetilde{\beta} = \alpha_0 + E,$$

where α_0 is the main term in $\mathcal{E}(\hat{\rho}^{-2}, \Lambda^{0,1}T^*\hat{X})$, and E is an error term in $\mathcal{E}(\hat{\rho}^0, \Lambda^{0,1}T^*\hat{X})$. Recall our goal was to compute $\Pi f = \Pi \overline{\partial}_{b,1}^*(\overline{\partial}_{b,1}\widetilde{\beta})$. But by Lemma 9.1,

$$\Pi \overline{\partial}_{b,1}^* E = 0.$$

Thus it suffices to compute $\Pi \overline{\partial}_{b,1}^* \alpha_0$. Now define, for $\varepsilon > 0$,

$$\alpha_{\varepsilon} := -2\pi i \chi v^{-1} \overline{\psi}^{-2} (\overline{\psi} + \varepsilon)^{-1} \hat{Z}^{\bar{1}}.$$

Then

$$\alpha_{\varepsilon} \in \mathcal{E}(\hat{\rho}^0, \Lambda^{0,1} T^* \hat{X}) \tag{9.1}$$

for all $\varepsilon > 0$, since by Theorem 4.4, Re $\psi \ge 0$, which implies $\overline{\psi} + \varepsilon \ne 0$ on X. By (1.31), there exists some $s \in \mathcal{E}(\hat{\rho}^3)$ such that

$$\bar{\partial}_{b,1}^*(h\hat{Z}^{\bar{1}}) = (-\hat{Z}_1 - \frac{\hat{Z}_1 v}{v} + s)\hbar$$

for all $h \in C^{\infty}(X)$. But $\hat{Z}_1 \overline{\psi} = 0$ on X. Thus

$$\overline{\partial}_{b,1}^* \alpha_0 = -2\pi i s \chi v^{-1} \overline{\psi}^{-3} + 2\pi i (\hat{Z}_1 \chi) v^{-1} \overline{\psi}^{-3} \quad \text{on } X.$$

Similarly,

$$\overline{\partial}_{b,1}^* \alpha_{\varepsilon} = -2\pi i s \chi v^{-1} \overline{\psi}^{-2} (\overline{\psi} + \varepsilon)^{-1} + 2\pi i (\hat{Z}_1 \chi) v^{-1} \overline{\psi}^{-2} (\overline{\psi} + \varepsilon)^{-1} \quad \text{on } X$$

It follows that

$$\overline{\partial}_{b,1}^* \alpha_{\varepsilon} \to \overline{\partial}_{b,1}^* \alpha_0 \quad \text{in } \mathcal{E}(\hat{\rho}^{-1}).$$

From (9.1) and Lemma 9.1, we then have

$$\Pi \overline{\partial}_{b,1}^* \alpha_0 = \lim_{\varepsilon \to 0} \Pi \overline{\partial}_{b,1}^* \alpha_\varepsilon = 0.$$

This completes the proof of Theorem 1.5.

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Appendix A. The Green's function of the conformal Laplacian

Assume, as in Section 1.1, that \hat{X} is a compact pseudohermitian 3-manifold of positive Tanaka–Webster class. In order to apply our results to the positive *p*-mass theorem in [5], one needs to check that the Green function G_p for the conformal sublaplacian $L_b := -4 \Delta_b + R$ at *p* satisfies (1.1) under the assumption that the Tanaka–Webster curvature *R* is positive on *X*. This can be done by using an argument similar to the one in Theorem 2.4. We recall that Δ_b denotes the sublaplacian on X. It was shown in Section 5 of [5] that G_p has the form: $G_p = \frac{1}{2\pi\hat{\sigma}^2} + \omega$, where $\omega \in C^1(\hat{X})$ and ω satisfies the equation

$$L_b \omega = \tilde{g}, \quad \tilde{g} \in \mathcal{E}(\hat{\rho}^0). \tag{A.1}$$

It is obvious that G_p satisfies (1.1) if $\omega - \omega(p) \in \mathcal{E}(\hat{\rho})$. We are going to prove that $\omega - \omega(p) \in \mathcal{E}(\hat{\rho})$.

First we extend L_b to

$$L_b: \operatorname{Dom}(L_b) \subset L^2(\hat{m}) \to L^2(\hat{m})$$

in the standard way. Note that $L_b = -4\triangle_b + R$ is subelliptic, self-adjoint and $-4\triangle_b + R$ has L^2 closed range. Since the conformal sublaplacian is conformally invariant, and the Tanaka–Webster curvature R can be made positive on X via a conformal change of the contact form, it follows that L_b : Dom $(L_b) \subset L^2(\hat{m}) \to L^2(\hat{m})$ is injective. Let $H: L^2(\hat{m}) \to \text{Dom}(L_b)$ be the inverse of L_b . We have

$$L_b H = I \quad \text{on } L^2(\hat{m}),$$

$$HL_b = I \quad \text{on Dom} (L_b). \tag{A.2}$$

We can repeat the L^2 estimates of Kohn and show that for every $k \in \mathbb{N}_0$, there is a constant $c_k > 0$ such that

$$\left\|\hat{\nabla}_{b}^{k+2}Hu\right\|_{\hat{m}} \le c_{k} \left\|\hat{\nabla}_{b}^{k}u\right\|_{\hat{m}}, \quad \forall u \in C^{\infty}(\hat{X}),$$
(A.3)

and the distribution kernel H(x, y) of H is C^{∞} away from the diagonal.

Now, we claim that H is a smoothing operator of order 2. Let B(x, r) be a small ball and let (x_1, x_2, x_3) be local coordinates on B(x, r). We first observe that for any smooth function f with support in B(x, r), we have

$$\begin{split} \|f(x)\|_{L^{\infty}(B(x,r))} &\leq c_0 \int\limits_{B(x,2r)} \left| \frac{\partial^3 f}{\partial x_1 \partial x_2 \partial x_3}(x) \right| \hat{m}(x) \\ &\leq c_1 \int\limits_{B(x,2r)} \left| \hat{\nabla}_b^4 f(x) \right| \hat{m}(x), \end{split}$$
(A.4)

where $c_0 > 0$ and $c_1 > 0$ are constants independent of f and r. Let ϕ be a normalized bump function in the ball B(x, r) and let $k \in \mathbb{N}_0$. From (A.3) and (A.4), we have

$$\begin{split} \left\| \hat{\nabla}_{b}^{k} H \phi \right\|_{L^{\infty}(B(x,r))} &\leq c_{1} \sum_{j=0}^{4} r^{j-4} \int_{B(x,2r)} \left| (\hat{\nabla}_{b}^{k+j} H \phi)(x) \right| \hat{m}(x) \\ &\leq c_{2} \sum_{j=0}^{4} r^{j-2} \left\| \hat{\nabla}_{b}^{k+j} H \phi \right\|_{\hat{m}} \leq \tilde{c}_{k} \sum_{j=0}^{4} r^{j-2} \left\| \hat{\nabla}_{b}^{k+j-2} \phi \right\|_{\hat{m}} \\ &\leq \hat{c}_{k} r^{2-k}, \end{split}$$
(A.5)

where $c_1 > 0$, $c_2 > 0$, $\tilde{c}_k > 0$ and $\hat{c}_k > 0$ are constants independent of r, ϕ and x. Thus, H satisfies the cancellation property for a smoothing operator of order 2. From (A.5) and (A.3), we can repeat the methods as in Christ [6,7] and Koenig [15] and conclude that for all multi-indices α_1 , α_2 , we have,

$$|(\nabla_b^{\alpha_1})_x(\nabla_b^{\alpha_2})_y H(x,y)| \le C_\alpha \vartheta(x,y)^{-2-|\alpha|}, \forall (x,y) \in \hat{X} \times \hat{X}, \quad x \neq y,$$
(A.6)

where $C_{\alpha} > 0$ is a constant. This shows that H is a smoothing operator of order 2.

Now, we are ready to prove that $\omega - \omega(p) \in \mathcal{E}(\hat{\rho})$. From (A.1) and (A.2), we have $\omega = H\tilde{g}$. Fix $k \in \mathbb{N}_0$ and fix a point $x_0 \neq p$, x_0 is in some small neighborhood W of p. Let $r = \frac{1}{4}\vartheta(x_0, p)$, and η be a normalized bump function supported in $B(x_0, r)$, with $\eta = 1$ on $B(x_0, r/2)$. Then,

$$|\hat{\nabla}_{b}^{k+1}\omega(x_{0})| = |\hat{\nabla}_{b}^{k+1}H\tilde{g}(x_{0})| \le |\hat{\nabla}_{b}^{k+1}H(\eta\tilde{g})(x_{0})| + |\hat{\nabla}_{b}^{k+1}H((1-\eta)\tilde{g})(x_{0})|$$

and $\eta \tilde{g}(x)$ is a normalized bump function on $B(x_0, r)$. So by (A.5), we see that

$$|\hat{\nabla}_b^{k+1} H(\eta \widetilde{g})(x_0)| \le C_k r^{1-k}, \tag{A.7}$$

where $C_k > 0$ is a constant independent of x_0 and r. By using (A.6), $\hat{\nabla}_b^{k+1} H((1-\eta)\tilde{g})(x_0)$ can be estimated by writing out the integral directly:

$$\hat{\nabla}_{b}^{k+1}H((1-\eta)\tilde{g})(x_{0}) = \int (\hat{\nabla}_{b}^{k+1}H)(x_{0},y)(1-\eta)(y)\tilde{g}(y)\hat{m}(y) + \int (\hat{\nabla}_{b}^{k+1}H)(x_{0},y)(1-\eta)(y)\hat{m}(y)\hat{m}(y)\hat{m}(y)\hat{m}(y) + \int (\hat{\nabla}_{b}^{k+1}H)(x_{0},y)(1-\eta)(y)\hat{m}(y)\hat$$

this integral is dominated by

$$D_k \int_{\vartheta(y,x_0) \ge \frac{1}{2}r} \vartheta(y,x_0)^{-3-k} \hat{m}(y) \le E_k r^{-k} \int_{\hat{X}} \vartheta(y,x_0)^{-3} \hat{m}(y) \le F_k r^{-k}, \quad (A.8)$$

where $D_k > 0$, $E_k > 0$ and $F_k > 0$ are constants independent of the point x_0 and r. From (A.7) and (A.8), we conclude that $\hat{\nabla}_b \omega \in \mathcal{E}(\hat{\rho}^0)$ and hence $\omega - \omega(p) \in \mathcal{E}(\hat{\rho})$.

Appendix B. Subelliptic estimates for $\overline{\partial}_b$

In this appendix we present a proof of Propositions 2.5 and 2.6. As is well-known, the crux of the matter is to prove a normalized subelliptic estimate on a unit cube, and rescale to a ball of radius r. It is this normalized subelliptic estimate we will focus on below.

Suppose on \mathbb{R}^3 , $T = \frac{\partial}{\partial x_3}$, and on the cube $Q_2 := (-2, 2)^3$, there is a (complex) vector field Z such that $[Z, \overline{Z}] = -iT + bZ + \overline{b}\overline{Z}$ and $[Z, T] = cZ + d\overline{Z} + eT$. Fix a sequence of positive numbers c_k . The only assumptions we make on b, c, d and e are that they are C^{∞} on Q_2 , and that their C^k norms are bounded by c_k for all k.

We also need to assume the following condition on Z: Write $Z = \sum_{i=1}^{3} A_i(x) \frac{\partial}{\partial x_i}$ on Q_2 . Then the only assumption we make on the A_i 's is that $|A_i(x)| \leq 1$.

Note that Z and \overline{Z} are only defined on Q_2 . They will never hit any function that is not supported in Q_2 , whereas T could hit a function that is defined on all of \mathbb{R}^3 .

Suppose also that we have a smooth contact form θ on Q_2 , so that $\theta(T) = 1$, and $\theta(Z) = \theta(\overline{Z}) = 0$ on Q_2 . One then has a measure $\theta \wedge d\theta$ on Q_2 . We also assume that the formal adjoint of Z with respect to $L^2(\theta \wedge d\theta)$ on Q_2 is given by $-\overline{Z} + a$ for some C^{∞} function a, where again the only assumption on a is that its C^k norm is bounded by c_k for all k. Similarly for the adjoint of \overline{Z} . We also assume that on Q_2 , $\theta \wedge d\theta = \rho^2 dx$ where $dx = dx_1 dx_2 dx_3$ is the Lebesgue measure on \mathbb{R}^3 and ρ is a positive smooth function on Q_2 . The only assumptions on ρ are that $c_0^{-1} \leq \rho \leq 1$ and that its C^1 norm is bounded by c_1 . There are no other assumptions on θ .

We shall also fix two functions η , $\tilde{\eta}$ such that they are C_c^{∞} with support in Q_2 , identically equal to 1 on $Q_1 := (-1, 1)^3$, and $\tilde{\eta} \equiv 1$ on a neighborhood of the support of η .

Write $\nabla_b u$ for $(Zu, \overline{Z}u)$. We claim the following proposition:

Proposition B.1. For all functions $u \in C^{\infty}(Q_2)$ and $k \ge 1$, we have

$$\|\nabla_b^k(\eta u)\| \le C_k \left(\|\nabla_b^{k-1}\overline{Z}(\tilde{\eta} u)\| + \|\tilde{\eta} u\| + \|\tilde{\eta} v\| \right)$$

where v is any solution to the equation Zv = u on Q_2 , all norms are $L^2(\theta \wedge d\theta)$ norms, and C_k depends only on the chosen sequence c_k and on η , $\tilde{\eta}$ (but not otherwise on the vector fields, the coefficients a, b, c, d, e, A_i or θ).

We remark that if we have $Zv + \alpha v = u$ instead of Zv = u, where α is a fixed C^{∞} function on Q_2 , then the above theorem still holds. See the end of this section for a discussion about that. Propositions 2.5 and 2.6 now follow easily by a well-known rescaling procedure. We omit the details.

To prove Proposition B.1, recall on \mathbb{R}^3 we have the Lebesgue measure $dx = dx_1 dx_2 dx_3$, and there is the Fourier transform defined by $\hat{u}(\xi) = \int u(x) e^{-2\pi i x \xi} dx$. Let Ψ^+ be a smooth function of ξ , such that it is (Euclidean) homogeneous of degree 0 outside the unit ball $\{|\xi| \leq 1\}$, equal to 1 on $\{\xi_3 > 2\varepsilon_0(|\xi_1| + |\xi_2|)\}$ there, and equal to 0 on both $\{\xi_3 < \varepsilon_0(|\xi_1| + |\xi_2|)\} \cap \{|\xi| > 1\}$ and $\{|\xi| \leq \frac{1}{2}\}$. Here ε_0 is a small positive absolute constant to be chosen. Let $\Psi^-(\xi) = \Psi^+(\xi_1, \xi_2, -\xi_3)$, and $\Psi^0 = 1 - \Psi^+ - \Psi^-$. Let $\Lambda^+, \Lambda^0, \Lambda^-$ be the Fourier multipliers corresponding to Ψ^+, Ψ^0, Ψ^- respectively. For instance, if U is a function on \mathbb{R}^3 , then $\widehat{\Lambda^+U}(\xi) := \Psi^+(\xi)\widehat{U}(\xi)$. Then $\Lambda^+ + \Lambda^0 + \Lambda^-$ is the identity operator on $L^2(dx)$.

We shall also fix a sequence of C^{∞} functions $\Psi_0^+, \Psi_1^+, \Psi_2^+, \ldots$ of ξ , with $\Psi_0^+ = \Psi^+$, such that each Ψ_k^+ is (Euclidean) homogeneous of degree 0 outside the unit ball $\{|\xi| \leq 1\}$, equal to 1 on a neighborhood of the support of Ψ_{k-1}^+ , and equal to 0 on both $\{\xi_3 < (1/2)\varepsilon_0(|\xi_1|+|\xi_2|)\} \cap \{|\xi| > 1\}$ and $\{|\xi| \leq \frac{1}{4}\}$. We shall denote the corresponding Fourier multipliers Λ_k^+ . Note $\Lambda_k^+ \Lambda_{k-1}^+ = \Lambda_{k-1}^+$ for all $k \geq 1$.

In addition, we fix a Fourier multiplier operator $\tilde{\Lambda}^-$, with symbol $\tilde{\Psi}^-$ that is supported on $\{\xi_3 < 0\}$, such that $\tilde{\Lambda}^- \Lambda^- = \Lambda^-$.

Finally, we fix a sequence of C_c^{∞} functions $\eta_0, \eta_1, \eta_2, \ldots$, with $\eta_0 = \eta$, such that each η_k has support in $Q_2, \eta_{k+1} \equiv 1$ on a neighborhood of the support of η_k for all k, and $\tilde{\eta} \equiv 1$ on the support of η_k for all k.

Now suppose we are given $u \in C^{\infty}(Q_2)$. We write u^+ for $\eta_1 \Lambda^+(\eta u)$, and similarly u^0 and u^- . Note that then $\eta u = u^+ + u^0 + u^-$.

B.1. Estimate for u^-

First we prove that for any $k \ge 1$,

$$\|\nabla_b^k u^-\| \le C_k \left(\|\nabla_b^{k-1} \overline{Z} u^-\| + \|\eta u\| \right).$$
(B.1)

A useful lemma is the following:

Lemma B.2. For every $k \ge 1$, there exist Euclidean pseudodifferential operators S_{-1} and S_{-k} , smoothing of orders 1 and k respectively, so that

$$u^{-} = \tilde{\Lambda}^{-} u^{-} + S_{-1} u^{-} + S_{-k} (\eta u).$$

Here a Euclidean pseudodifferential operator is said to be smoothing of order k, if its symbol is in the Hörmander class $S_{1,0}^{-k}$.

Proof. To see this, write

$$u^{-} = \eta_{1}\Lambda^{-}\eta u$$

= $\eta_{1}\tilde{\Lambda}^{-}\eta_{2}\Lambda^{-}\eta u + \eta_{1}\tilde{\Lambda}^{-}(1-\eta_{2})\Lambda^{-}\eta u$
= $\tilde{\Lambda}^{-}\eta_{1}\Lambda^{-}\eta u + [\eta_{1},\tilde{\Lambda}^{-}]\eta_{2}\Lambda^{-}\eta u + \eta_{1}\tilde{\Lambda}^{-}(1-\eta_{2})\Lambda^{-}\eta u.$ (B.2)

The last term here is $S_{-k}(\eta u)$ for some (Euclidean) pseudodifferential operator that is smoothing of order k, because one can pick some C_c^{∞} function ζ such that $\zeta \equiv 1$ on
the support of η and $\eta_1 \equiv 1$ on the support of ζ ; this is possible because $\eta_1 \equiv 1$ on a neighborhood of the support of η . Then writing ηu as $\zeta \eta u$, and commuting the ζ past Λ^- to hit $(1 - \eta_2)$, the last term above is just

$$\eta_1 \tilde{\Lambda}^- (1 - \eta_2) [[[\Lambda^-, \zeta], \zeta], \dots, \zeta](\eta u) = S_{-k}(\eta u).$$

With the same choice of ζ , the second term in (B.2) can be written as

$$[\eta_1, \tilde{\Lambda}^-]\eta_2\eta_1\Lambda^-\eta u + [\eta_1, \tilde{\Lambda}^-]\eta_2(1-\eta_1)\Lambda^-\eta u = [\eta_1, \tilde{\Lambda}^-]u^- + [\eta_1, \tilde{\Lambda}^-]\eta_2(1-\eta_1)\Lambda^-\eta u,$$

and by the same argument above, the last term here is a pseudodifferential operator of order -k acting on ηu . It follows that this second term in (B.2) is of the form $S_{-1}u^- + S_{-k}(\eta u)$. Finally, the first term in (B.2) is just $\tilde{\Lambda}^- u^-$. This completes our proof of this lemma. \Box

Also, to prove (B.1), it suffices to prove that

$$\|\nabla_b^k u^-\| \le C_k \left(\sum_{l=0}^{k-1} \|\nabla_b^l \overline{Z} u^-\| + \|\eta u\|\right),$$
 (B.3)

in view of the following interpolation inequality:

Lemma B.3. If $k \ge 1$, then

$$\|\nabla_b^l \overline{Z}u\|^2 \le C_k \left(\|\nabla_b^{k-1} \overline{Z}u\|^2 + \|u\|^2\right)$$

for all $0 \leq l \leq k-1$, for any function u that is smooth and compactly supported in Q_2 .

Proof. One proves, by induction on l beginning at l = 0, that for any $\varepsilon > 0$, there exists $C_{k,\varepsilon}$ such that

$$\|\nabla_b^l \overline{Z}u\|^2 \le \varepsilon \sum_{j=0}^{k-1} \|\nabla_b^j \overline{Z}u\|^2 + C_{k,\varepsilon} \|u\|^2$$

for all $0 \leq l \leq k - 2$. The key is that

$$\begin{aligned} (\nabla_b^l \overline{Z}u, \nabla_b^l \overline{Z}u) &= (\nabla_b^{l-1} \overline{Z}u, \nabla_b^{l+1} \overline{Z}u) + O(\|\nabla_b^{l-1}u\| \|\nabla_b^l u\|) \\ &\leq \varepsilon (\|\nabla_b^{l+1} \overline{Z}u\|^2 + \|\nabla_b^l \overline{Z}u\|^2) + C_{l,\varepsilon} \|\nabla_b^{l-1} \overline{Z}u\|^2. \end{aligned}$$

Once this is established, the lemma follows easily by summing over l. \Box

Now, to prove (B.3), we proceed by induction on k.

When k = 1, it suffices to bound $||Zu^-||^2$. Denote (\cdot, \cdot) the inner product in $L^2(\theta \wedge d\theta)$. Then

$$||Zu^{-}||^{2} = -(\overline{Z}Zu^{-}, u^{-}) + O(||u^{-}|| ||\nabla_{b}u^{-}||)$$

$$= -(Z\overline{Z}u^{-}, u^{-}) - (iTu^{-}, u^{-}) + O(||u^{-}|| ||\nabla_{b}u^{-}||)$$

$$= ||\overline{Z}u^{-}||^{2} - (iTu^{-}, u^{-}) + O(||u^{-}|| ||\nabla_{b}u^{-}||).$$

Now $(iTu^-, u^-) = \langle \rho iTu^-, \rho u^- \rangle$, where $\langle \cdot, \cdot \rangle$ is the L^2 inner product with respect to the Lebesgue measure dx. Hence

$$(iTu^{-}, u^{-}) = \langle iT\rho u^{-}, \rho u^{-} \rangle + O(||u^{-}||^{2}).$$

Also, by the above lemma,

$$\rho u^{-} = \tilde{\Lambda}^{-} \rho u^{-} + [\rho, \tilde{\Lambda}^{-}] u^{-} + \rho S_{-1} u^{-} + \rho S_{-1} (\eta u)$$

= $\tilde{\Lambda}^{-} \rho u^{-} + S_{-1} u^{-} + S_{-1} (\eta u).$

 So

$$\begin{split} (iTu^-, u^-) &= \langle iT\tilde{\Lambda}^-\rho u^-, \tilde{\Lambda}^-\rho u^- \rangle + \langle iT\tilde{\Lambda}^-\rho u^-, S_{-1}u^- \rangle \\ &+ \langle iT\tilde{\Lambda}^-\rho u^-, S_{-1}(\eta u) \rangle + O(\|u^-\|^2) + O(\|\eta u\|^2). \end{split}$$

But the second and third terms are $O(||u^-||^2) + O(||\eta u||^2)$ (one just needs to integrate by parts in T and let T fall on S_{-1}), and the first term is

$$\langle iT\tilde{\Lambda}^-\rho u^-, \tilde{\Lambda}^-\rho u^- \rangle = \int -2\pi\xi_3 |\tilde{\Psi}^-(\xi)|^2 |\widehat{\rho u^-}(\xi)|^2 d\xi$$

which is non-negative since $\xi_3 < 0$ on the support of $\tilde{\Psi}^-$. Hence altogether

$$||Zu^{-}||^{2} \leq C||\overline{Z}u^{-}||^{2} + O(||u^{-}||^{2}) + O(||\eta u||^{2}) + O(||u^{-}||||\nabla_{b}u^{-}||),$$

and using $||u^-|| ||\nabla_b u^-|| \le \delta ||\nabla_b u^-||^2 + \delta^{-1} ||u^-||^2$ and $||u^-||^2 \le C ||\eta u||^2$, we get

$$\|Zu^-\|^2 \le C(\|\overline{Z}u^-\|^2 + \|\eta u\|^2)$$

as desired.

Next, suppose (B.3) has been proved for k-1 for some $k \ge 2$. We prove the same estimate for k. To do so, we first prove that for all $0 \le m \le \lfloor \frac{k}{2} \rfloor$ and all $\varepsilon > 0$, there exists C_{ε} such that

$$\|T^{m}\nabla_{b}^{k-2m}u^{-}\|^{2} \leq \varepsilon \|\nabla_{b}^{k}u^{-}\|^{2} + C_{\varepsilon} \left(\|\nabla_{b}^{k-1}\overline{Z}u^{-}\|^{2} + \sum_{l=0}^{k-1} \|\nabla_{b}^{l}u^{-}\|^{2} + \|\eta u\|^{2}\right).$$
(B.4)

In fact the desired inequality (B.1) for k follows readily from the above inequality when m = 0.

To prove (B.4), we proceed in two steps. First, we prove that for all $0 \le m \le \lfloor \frac{k}{2} \rfloor$ and all $\varepsilon > 0$, there exists C_{ε} such that

$$\|T^{m}Z^{k-2m}u^{-}\|^{2} \leq \varepsilon \|\nabla_{b}^{k}u^{-}\|^{2} + C_{\varepsilon} \left(\|\nabla_{b}^{k-1}\overline{Z}u^{-}\|^{2} + \sum_{l=0}^{k-1} \|\nabla_{b}^{l}u^{-}\|^{2} + \|\eta u\|^{2} \right).$$
(B.5)

Next, we prove by induction on m, beginning from m that is as large as possible, that (B.4) holds.

In the first step, there are two cases: either k - 2m = 0 (which occurs only when k is even), or $k - 2m \ge 1$.

In the first case, we need to estimate $||T^m u^-||^2$. Now $m \ge 1$, and

$$\begin{split} \|T^{m}u^{-}\|^{2} &= (T^{m}u^{-}, i(Z\overline{Z} - \overline{Z}Z)T^{m-1}u^{-}) + O(\|T^{m}u^{-}\|\|\nabla_{b}^{k-1}u^{-}\|) \\ &= -(T^{m}u^{-}, i\overline{Z}ZT^{m-1}u^{-}) + (T^{m}u^{-}, iZT^{m-1}\overline{Z}u^{-}) + O(\|T^{m}u^{-}\|\sum_{l=0}^{k-1}\|\nabla_{b}^{l}u^{-}\|) \\ &= -(T^{m}u^{-}, i\overline{Z}ZT^{m-1}u^{-}) + O(\|T^{m}u^{-}\|\|\nabla_{b}^{k-1}\overline{Z}u^{-}\|) + O(\|T^{m}u^{-}\|\sum_{l=0}^{k-1}\|\nabla_{b}^{l}u^{-}\|). \end{split}$$

But

$$\begin{split} (T^{m}u^{-}, i\overline{Z}ZT^{m-1}u^{-}) \\ &= (iT(ZT^{m-1}u^{-}), ZT^{m-1}u^{-}) + O(\|T^{m}u^{-}\|\|\nabla_{b}^{k-1}u^{-}\|) + O(\|\nabla_{b}^{k-1}u^{-}\|^{2}) \\ &= \langle iT(\rho ZT^{m-1}u^{-}), \rho ZT^{m-1}u^{-} \rangle + O(\|T^{m}u^{-}\|\|\nabla_{b}^{k-1}u^{-}\|) + O(\|\nabla_{b}^{k-1}u^{-}\|^{2}). \end{split}$$

Also by Lemma B.2,

$$\begin{split} \rho Z T^{m-1} u^{-} \\ &= \tilde{\Lambda}^{-} (\rho Z T^{m-1} u^{-}) + [\rho Z T^{m-1}, \tilde{\Lambda}^{-}] u^{-} + \rho Z T^{m-1} S_{-1}(u^{-}) + \rho Z T^{m-1} S_{-k}(\eta u) \\ &= \tilde{\Lambda}^{-} (\rho Z T^{m-1} u^{-}) + S_{-1} \sum_{l=0}^{k-1} \nabla_{b}^{l} u^{-} + S_{-1}(\eta u). \end{split}$$

Hence

$$\begin{split} \langle iT(\rho ZT^{m-1}u^{-}), \rho ZT^{m-1}u^{-} \rangle \\ &= \langle iT\tilde{\Lambda}^{-}(\rho ZT^{m-1}u^{-}), \tilde{\Lambda}^{-}(\rho ZT^{m-1}u^{-}) \rangle + \langle iT\tilde{\Lambda}^{-}(\rho ZT^{m-1}u^{-}), S_{-1}\sum_{l=0}^{k-1} \nabla_{b}^{l}u^{-} \rangle \\ &+ \langle iT\tilde{\Lambda}^{-}(\rho ZT^{m-1}u^{-}), S_{-1}(\eta u) \rangle + O(\sum_{l=0}^{k-1} \|\nabla_{b}^{l}u^{-}\|^{2}) + O(\|\eta u\|^{2}), \end{split}$$

where the first term is non-negative, and the second and third terms are $O(\sum_{l=0}^{k-1} \|\nabla_b^l u^-\|^2) + O(\|\eta u\|^2)$ after integrating by parts in T. Altogether, we get

$$\|T^{m}u^{-}\|^{2} \leq \varepsilon \|\nabla_{b}^{k}u^{-}\|^{2} + C_{\varepsilon} \left(\|\nabla_{b}^{k-1}\overline{Z}u^{-}\|^{2} + \sum_{l=0}^{k-1} \|\nabla_{b}^{l}u^{-}\|^{2} + \|\eta u\|^{2} \right)$$

as desired.

Next, in the second case, we need to estimate $||T^m Z^{k-2m} u^-||^2$ when $k-2m \ge 1$. The strategy is the same as the one when we dealt with the case when k = 1 and m = 0. One observes that

$$\begin{split} \|T^m Z^{k-2m} u^-\|^2 \\ &= -(T^m \overline{Z} Z^{k-2m} u^-, T^m Z^{k-2m-1} u^-) + O(\sum_{l=0}^{k-1} \|\nabla_b^l u^-\| \|\nabla_b^k u^-\|) + O(\sum_{l=0}^{k-1} \|\nabla_b^l u^-\|^2) \\ &= -(T^m Z^{k-2m} \overline{Z} u^-, T^m Z^{k-2m-1} u^-) \\ &- (k-2m)(iT(T^m Z^{k-2m-1} u^-), T^m Z^{k-2m-1} u^-) \\ &+ O(\sum_{l=0}^{k-1} \|\nabla_b^l u^-\| \|\nabla_b^k u^-\|) + O(\sum_{l=0}^{k-1} \|\nabla_b^l u^-\|^2) \\ &= \|T^m Z^{k-2m-1} \overline{Z} u^-\|^2 - [2(k-2m)-1]\langle iT(\rho T^m Z^{k-2m-1} u^-), \rho T^m Z^{k-2m-1} u^-\rangle \\ &+ O(\sum_{l=0}^{k-1} \|\nabla_b^l u^-\| \|\nabla_b^k u^-\|) + O(\sum_{l=0}^{k-1} \|\nabla_b^l u^-\|^2). \end{split}$$

Now [2(k-2m)-1] > 0, and by Lemma B.2,

$$\begin{split} \rho T^m Z^{k-2m-1} u^- &= \tilde{\Lambda}^- (\rho T^m Z^{k-2m-1} u^-) + [\rho T^m Z^{k-2m-1}, \tilde{\Lambda}^-] u^- \\ &+ \rho T^m Z^{k-2m-1} S_{-1}(u^-) + \rho T^m Z^{k-2m-1} S_{-k}(\eta u) \\ &= \tilde{\Lambda}^- (\rho T^m Z^{k-2m-1} u^-) + S_{-1} \sum_{l=0}^{k-1} \nabla_b^l u^- + S_{-1}(\eta u). \end{split}$$

Hence

$$\begin{split} \langle iT(\rho T^m Z^{k-2m-1} u^-), \rho T^m Z^{k-2m-1} u^- \rangle \\ &= \langle iT\tilde{\Lambda}^-(\rho T^m Z^{k-2m-1} u^-), \tilde{\Lambda}^-(\rho T^m Z^{k-2m-1} u^-) \rangle \\ &+ \langle iT\tilde{\Lambda}^-(\rho T^m Z^{k-2m-1} u^-), S_{-1} \sum_{l=0}^{k-1} \nabla_b^l u^- \rangle \\ &+ \langle iT\tilde{\Lambda}^-(\rho T^m Z^{k-2m-1} u^-), S_{-1}(\eta u) \rangle + O(\sum_{l=0}^{k-1} \|\nabla_b^l u^-\|^2) + O(\|\eta u\|^2), \end{split}$$

where the first term is non-negative, and the second and third terms are $O(\sum_{l=0}^{k-1} \|\nabla_b^l u^-\|^2) + O(\|\eta u\|^2)$ after integrating by parts in T. Altogether,

$$\|T^m Z^{k-2m} u^-\|^2 \le \varepsilon \|\nabla_b^k u^-\|^2 + C_{\varepsilon} \left(\|\nabla_b^{k-1} \overline{Z} u^-\|^2 + \sum_{l=0}^{k-1} \|\nabla_b^l u^-\|^2 + \|\eta u\|^2 \right)$$

as desired. This finishes our first step in proving (B.4).

Now to complete the proof of (B.4), we proceed by induction on m, beginning with m that is as big as possible. In that case k - 2m is either 0 or 1. Both cases follow right away by what we have proved above in the first step. Now we prove the inequality (B.4) for m, assuming the inequality has been proved for all strictly bigger m's. Then we need to estimate $||T^m \nabla_b^{k-2m} u^-||^2$. Consider $T^m \nabla_b^{k-2m} u^-$. If all the ∇_b 's are Z, then this follows again from what we have proved above. If one of the ∇_b 's to get $T^m \nabla_b^{k-2m-1} \overline{Z} u^-$, up to an error that either has fewer ∇_b derivatives, or an error of the form $T^{m+1} \nabla_b^{k-2m-2} u^-$. For example,

$$\begin{split} \|T^m \overline{Z} \nabla_b^{k-2m-1} u^-\|^2 \\ &= \|T^m \nabla_b^{k-2m-1} \overline{Z} u^-\|^2 + O(\|T^{m+1} \nabla_b^{k-2m-2} u^-\|^2) + O(\sum_{l=0}^{k-1} \|\nabla_b^l u^-\|^2). \end{split}$$

The first error term can then be estimated by our induction hypothesis on m. Hence

$$\|T^m \overline{Z} \nabla_b^{k-2m-1} u^-\|^2 \le \varepsilon \|\nabla_b^k u^-\|^2 + C_{\varepsilon} \left(\|\nabla_b^{k-1} \overline{Z} u^-\|^2 + \sum_{l=0}^{k-1} \|\nabla_b^l u^-\|^2 + \|\eta u\|^2 \right).$$

This completes the proof of (B.4), and thus the proof of (B.1).

B.2. Estimate for u^0

Next we prove that for any $k \ge 1$ and any $\varepsilon > 0$,

$$\|\nabla_{b}^{k} u^{0}\| \leq \|\nabla_{b}^{k-1} \overline{Z} u^{0}\|^{2} + \varepsilon \|\nabla_{b}^{k} (\eta u)\|^{2} + C_{k,\varepsilon} \sum_{l=0}^{k-1} \|\nabla_{b}^{l} (\eta u)\|.$$
(B.6)

We proceed by induction on k exactly as before.

When k = 1, we only need to estimate $|(iTu^0, u^0)|$. But

$$|(iTu^0, u^0)| \le \varepsilon ||Tu^0||^2 + \varepsilon^{-1} ||\eta u||^2,$$

and

$$||Tu^0||^2 \le ||T\Lambda^0(\rho\eta u)||^2 + C||\eta u||^2.$$

Taking Fourier transform,

$$\begin{split} \|T\Lambda^{0}(\rho\eta u)\|^{2} &\leq \int (1+2\varepsilon_{0}^{2}(|\xi_{1}|^{2}+|\xi_{2}|^{2}))|\widehat{\Lambda^{0}(\rho\eta u)}(\xi)|^{2}d\xi \\ &\leq 2\varepsilon_{0}^{2}\left(\left\|\frac{\partial}{\partial x_{1}}\Lambda^{0}(\rho\eta u)\right\|^{2}+\left\|\frac{\partial}{\partial x_{2}}\Lambda^{0}(\rho\eta u)\right\|^{2}\right)+\|\eta u\|^{2} \\ &\leq 2\varepsilon_{0}^{2}\left(\left\|\frac{\partial}{\partial x_{1}}\eta_{1}\Lambda^{0}(\rho\eta u)\right\|^{2}+\left\|\frac{\partial}{\partial x_{2}}\eta_{1}\Lambda^{0}(\rho\eta u)\right\|^{2}\right)+\|\eta u\|^{2}. \end{split}$$

(The last line follows since $\rho\eta u = \eta_1 \rho\eta u$ and one can commute the η_1 past Λ^0 to obtain a better error.)

Now the key observation is that on Q_2 , $\frac{\partial}{\partial x_1}$ and $\frac{\partial}{\partial x_2}$ can be written as linear combinations of Z, \overline{Z} and T with coefficients that are bounded by an absolute constant. In fact we only need to bound the coefficients of the inverse of the matrix whose first column is (A_1, A_2, A_3) , the second column is the conjugate of the first, and the third column is (0, 0, 1). From $\theta \wedge d\theta(Z, \overline{Z}, T) = d\theta(Z, \overline{Z}) = -\theta(Z, \overline{Z}) = \theta(iT) = i, \theta \wedge d\theta = \rho^2 dx$, and $\rho \leq 1$, we have $|dx(Z, \overline{Z}, T)| \geq 1$, i.e. the determinant of the matrix to be inverted is bounded below by 1. Together with the assumed bounds on the A_i 's, we obtain our key observation.

Hence, continuing from above,

$$\|T\Lambda^{0}(\rho\eta u)\|^{2} \leq C_{0}\varepsilon_{0}^{2}\left(\left\|\nabla_{b}\Lambda^{0}(\rho\eta u)\right\|^{2} + \left\|T\Lambda^{0}(\rho\eta u)\right\|^{2}\right) + \|\eta u\|^{2},$$

which implies

$$\|T\Lambda^{0}(\rho\eta u)\|^{2} \leq C \|\nabla_{b}\Lambda^{0}(\rho\eta u)\|^{2} + \|\eta u\|^{2}$$
$$\leq C(\|\nabla_{b}(\eta u)\|^{2} + \|\eta u\|^{2})$$

if ε_0 was chosen to be sufficiently small. One then completes the proof of the case k = 1 as before.

Next, we prove by induction on m the following for any $\varepsilon > 0$:

$$\|T^{m}\nabla_{b}^{k-2m}u^{0}\|^{2} \leq C_{\varepsilon} \left(\|\nabla_{b}^{k-1}\overline{Z}u^{0}\|^{2} + \sum_{l=0}^{k-1}\|\nabla_{b}^{l}(\eta u)\|^{2}\right) + \varepsilon\|\nabla_{b}^{k}(\eta u)\|^{2}.$$
(B.7)

First, suppose k - 2m = 0. Then we need only estimate $|(iTZT^{m-1}u^0, ZT^{m-1}u^0)|$. One certainly has

$$|(iTZT^{m-1}u^0, ZT^{m-1}u^0)| \le \varepsilon ||TZT^{m-1}u^0||^2 + \varepsilon^{-1} \sum_{l=0}^{k-1} ||\nabla_b^l(\eta u)||^2.$$

To estimate $||TZT^{m-1}u^0||^2$, we write $TZT^{m-1}u^0 = T\Lambda^0 ZT^{m-1}(\eta u) + T[ZT^{m-1}\eta_1, \Lambda^0](\eta u)$. The second term is bounded by

$$\sum_{l=0}^k \|\nabla_b^l(\eta u)\|$$

The first term satisfies

$$\|T\Lambda^0 Z T^{m-1}(\eta u)\|^2 \le \|\nabla_b^{k-1}(\eta u)\|^2 + C_0 \varepsilon_0^2 \|\nabla_b \Lambda^0 Z T^{m-1}(\eta u)\|^2 + C_0 \varepsilon_0^2 \|T\Lambda^0 Z T^{m-1}(\eta u)\|^2,$$

from which it follows that

$$\begin{aligned} \|T\Lambda^0 Z T^{m-1}(\eta u)\|^2 &\leq \|\nabla_b^{k-1}(\eta u)\|^2 + C \|\nabla_b \Lambda^0 Z T^{m-1}(\eta u)\|^2 \\ &\leq C \|\nabla_b^{k-1}(\eta u)\|^2 + C \|\nabla_b^k(\eta u)\|^2 \end{aligned}$$

by our choice of ε_0 . Hence

$$|(iTZT^{m-1}u^0, ZT^{m-1}u^0)| \le \varepsilon \|\nabla_b^k(\eta u)\|^2 + C_{\varepsilon} \sum_{l=0}^{k-1} \|\nabla_b^l(\eta u)\|^2$$

and one finishes the proof for the case k - 2m = 0 as before.

Next, when k - 2m = 1, we need to estimate $|(iTT^m u^0, T^m u^0)|$. One certainly has

$$|(iTT^{m}u^{0}, T^{m}u^{0})| \leq \varepsilon ||T^{m+1}u^{0}||^{2} + \varepsilon^{-1} \sum_{l=0}^{k-1} ||\nabla_{b}^{l}(\eta u)||^{2}.$$

To estimate $||T^{m+1}u^0||^2$, we write $T^{m+1}u^0 = T\Lambda^0 T^m(\eta u) + T[T^m\eta_1, \Lambda^0](\eta u)$. The second term is bounded by

$$\sum_{l=0}^k \|\nabla_b^l(\eta u)\|.$$

The first term satisfies

$$\|T\Lambda^0 T^m(\eta u)\|^2 \le \|\nabla_b^{k-1}(\eta u)\|^2 + C_0 \varepsilon_0^2 \|\nabla_b \Lambda^0 T^m(\eta u)\|^2 + C_0 \varepsilon_0^2 \|T\Lambda^0 T^m(\eta u)\|^2,$$

from which it follows that

$$\|T\Lambda^0 T^m(\eta u)\|^2 \le \|\nabla_b^{k-1}(\eta u)\|^2 + C\|\nabla_b\Lambda^0 T^m(\eta u)\|^2$$
$$\le C\|\nabla_b^{k-1}(\eta u)\|^2 + C\|\nabla_b^k(\eta u)\|^2.$$

Hence

$$|(iTT^{m}u^{0}, T^{m}u^{0})| \le \varepsilon \|\nabla_{b}^{k}(\eta u)\|^{2} + C_{\varepsilon} \sum_{l=0}^{k-1} \|\nabla_{b}^{l}(\eta u)\|^{2}$$

and one finishes the proof for the case k - 2m = 1 as before.

Now we prove (B.7) for m, assuming that the statement has been proved for all larger m's. We then estimate $||T^m \nabla_b^{k-2m} u^0||^2$. If one of the ∇_b is \overline{Z} , we proceed exactly as before and commute the \overline{Z} until it hits u^0 . This proves the desired estimate with the induction hypothesis on m. If now all ∇_b are Z's, then as before we only need to bound $|(iTT^m Z^{k-2m-1} u^0, T^m Z^{k-2m-1} u^0)|$. One certainly has

$$|(iTT^mZ^{k-2m-1}u^0, T^mZ^{k-2m-1}u^0)| \le \varepsilon \|T^{m+1}Z^{k-2m-1}u^0\|^2 + \varepsilon^{-1}\sum_{l=0}^{k-1} \|\nabla_b^l(\eta u)\|^2.$$

To estimate $||T^{m+1}Z^{k-2m-1}u^0||^2$, we write $T^{m+1}Z^{k-2m-1}u^0 = T\Lambda^0 T^m Z^{k-2m-1}(\eta u) + T[T^m Z^{k-2m-1}\eta_1, \Lambda^0](\eta u)$. The second term is bounded by

$$\sum_{l=0}^k \|\nabla_b^l(\eta u)\|.$$

The first term satisfies

$$\|T\Lambda^{0}T^{m}Z^{k-2m-1}(\eta u)\|^{2} \leq \|\nabla_{b}^{k-1}(\eta u)\|^{2} + C_{0}\varepsilon_{0}^{2}\|\nabla_{b}\Lambda^{0}T^{m}Z^{k-2m-1}(\eta u)\|^{2} + C_{0}\varepsilon_{0}^{2}\|T\Lambda^{0}T^{m}Z^{k-2m-1}(\eta u)\|^{2},$$

from which it follows that

$$\|T\Lambda^{0}T^{m}Z^{k-2m-1}(\eta u)\|^{2} \leq \|\nabla_{b}^{k-1}(\eta u)\|^{2} + C\|\nabla_{b}\Lambda^{0}T^{m}Z^{k-2m-1}(\eta u)\|^{2}$$
$$\leq C\|\nabla_{b}^{k-1}(\eta u)\|^{2} + C\|\nabla_{b}^{k}(\eta u)\|^{2}.$$

Hence

$$|(iTT^{m}Z^{k-2m-1}u^{0}, T^{m}Z^{k-2m-1}u^{0})| \le \varepsilon \|\nabla_{b}^{k}(\eta u)\|^{2} + C_{\varepsilon} \sum_{l=0}^{k-1} \|\nabla_{b}^{l}(\eta u)\|^{2}$$

and one finishes the proof for this case as before.

B.3. Estimate for u^+

Now we turn to estimate u^+ . Recall we introduced a sequence of cut-offs $\eta_0 = \eta, \eta_1, \eta_2, \ldots$, and a sequence of Fourier multipliers $\Lambda_0^+ = \Lambda^+, \Lambda_1^+, \Lambda_2^+, \ldots$, such that $\eta_k \eta_{k+1} = \eta_k$, and $\Lambda_k^+ \Lambda_{k+1}^+ = \Lambda_k^+$ for all $k \ge 0$. In fact the Fourier multiplier for Λ_{k+1}^+ is identically equal to 1 on a neighborhood of the support of that of Λ_k^+ . Also, the cut-off function $\tilde{\eta}$ dominates all the η_j 's, in the sense that $\eta_j \tilde{\eta} = \eta_j$ for all j. We wrote $u^+ = \eta_1 \Lambda^+ \eta u$, and v is any solution to Zv = u on Q_2 . The estimate we shall prove is

$$\|\nabla_b^k u^+\| \le C_k(\|\nabla_b^{k-1}\overline{Z}u^+\| + \sum_{l=0}^{k-1} \|\nabla_b^l(\eta_k u)\| + \|\eta_{k+1}v\|)$$
(B.8)

for all $k \geq 1$.

To prove this, the first observation is the following:

Lemma B.4. For all $k \geq 1$,

 $\|\nabla_{b}^{k}v^{+}\| \leq C_{k} \left(\|\nabla_{b}^{k-1}Zv^{+}\| + \|\eta v\| \right),$

where $v^+ = \eta_1 \Lambda^+ \eta v$.

The proof of this inequality is the same as that of (B.1), except that one reverses the role of Z and \overline{Z} , and replaces u^- by v^+ . It does not make use of the fact that v solves Zv = u. By the same token,

Lemma B.5. For all $k \ge 1$ and all $j \ge 1$,

$$\|\nabla_b^k(\eta_j \Lambda_{j-1}^+ \eta_{j-1} v)\| \le C_{j,k} \left(\|\nabla_b^{k-1} Z(\eta_j \Lambda_{j-1}^+ \eta_{j-1} v)\| + \|\eta_{j-1} v\| \right).$$

Another useful lemma is the following:

Lemma B.6. For any $k \ge 1$ and any $j \ge 1$, there exist pseudodifferential operators S_0 and S_{-k} , smoothing of order 0 and k respectively, such that

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$$\eta_j \Lambda_{j-1}^+ \eta_{j-1} u = Z(\eta_j \Lambda_j^+ \eta_{j-1} v) + S_0(\eta_{j+1} \Lambda_j^+ \eta_j v) + S_{-k}(\eta_j v).$$

In particular, when j = 1,

$$u^{+} = Zv^{+} + S_{0}(\eta_{2}\Lambda_{1}^{+}\eta_{1}v) + S_{-k}(\eta_{1}v).$$

Proof. For all $j \ge 1$,

$$\eta_{j}\Lambda_{j-1}^{+}\eta_{j-1}u = \eta_{j}\Lambda_{j-1}^{+}\eta_{j-1}Zv$$

$$= \eta_{j}\Lambda_{j-1}^{+}\eta_{j-1}Z(\eta_{j}v)$$

$$= \eta_{j}\Lambda_{j-1}^{+}\eta_{j-1}Z(\eta_{j+1}\Lambda_{j}^{+}\eta_{j}v) + \eta_{j}\Lambda_{j-1}^{+}\eta_{j-1}Z(\eta_{j+1}(1-\Lambda_{j}^{+})\eta_{j}v).$$
(B.9)

We shall argue that the second term on the last line is $S_{-k}(\eta_j v)$ for any $k \ge 1$.

Since the Fourier multiplier Ψ_j^+ of Λ_j^+ is identically 1 on a neighborhood of the support of Ψ_{j-1}^+ , there exists a Fourier multiplier $\Psi_{j,0}^+$ such that Ψ_j^+ is identically 1 on the support of $\Psi_{j,0}^+$, and $\Psi_{j,0}^+$ is identically 1 on the support of Ψ_{j-1}^+ . Writing $\Lambda_{j,0}^+$ for the Fourier multiplier operator corresponding to $\Psi_{j,0}^+$, we have

$$(1 - \Lambda_j^+) = (1 - \Lambda_{j,0}^+)(1 - \Lambda_j^+);$$

indeed $\Psi_j^+ \equiv 1$ on the support of $\Psi_{j,0}^+$ implies that $\Lambda_{j,0}^+(1-\Lambda_j^+) = 0$. Putting this back in the second term (B.9), and commuting $1 - \Lambda_{j,0}^+$ until it hits Λ_{j-1}^+ , we get

$$\eta_j \Lambda_{j-1}^+ (1 - \Lambda_{j,0}^+) \eta_{j-1} Z(\eta_{j+1} (1 - \Lambda_j^+) \eta_j v) + \eta_j \Lambda_{j-1}^+ [\eta_{j-1} Z \eta_{j+1}, 1 - \Lambda_{j,0}^+] (1 - \Lambda_j^+) \eta_j v).$$

The first term here is zero, since

$$\Lambda_{j-1}^{+}(1 - \Lambda_{j,0}^{+}) = 0;$$

the second term here is

$$\eta_j \Lambda_{j-1}^+ S_0 (1 - \Lambda_j^+) \eta_j v.$$

Again writing $(1 - \Lambda_j^+) = (1 - \Lambda_{j,0}^+)(1 - \Lambda_j^+)$ and commuting $1 - \Lambda_{j,0}^+$ until it hits Λ_{j-1}^+ , we get that this is

$$\eta_j \Lambda_{j-1}^+ S_{-1} (1 - \Lambda_j^+) \eta_j v.$$

Repeating this argument, it is clear that we can make this $\eta_j \Lambda_{j-1}^+ S_{-k} (1 - \Lambda_j^+) \eta_j v$ for any k, and this is thus $S_{-k}(\eta_j v)$.

Next, the first term in (B.9) is

$$\begin{split} \eta_{j}\Lambda_{j-1}^{+}\eta_{j-1}Z(\eta_{j+1}\Lambda_{j}^{+}\eta_{j}v) \\ &= Z\eta_{j}\Lambda_{j-1}^{+}\eta_{j-1}\eta_{j+1}\Lambda_{j}^{+}\eta_{j}v + [\eta_{j}\Lambda_{j-1}^{+}\eta_{j-1}, Z](\eta_{j+1}\Lambda_{j}^{+}\eta_{j}v) \\ &= Z\eta_{j}\Lambda_{j-1}^{+}\eta_{j-1}\Lambda_{j}^{+}\eta_{j}v + S_{0}(\eta_{j+1}\Lambda_{j}^{+}\eta_{j}v). \end{split}$$

By writing $\Lambda_j^+ = 1 - (1 - \Lambda_j^+)$, the first term in the last line is equal to

$$Z(\eta_{j}\Lambda_{j-1}^{+}\eta_{j-1}v) - Z\eta_{j}\Lambda_{j-1}^{+}\eta_{j-1}(1-\Lambda_{j}^{+})\eta_{j}v.$$

We only need to argue now that the second term in the last line is $S_{-k}(\eta_j v)$ for any k. But we only need to adopt the strategy above again: writing $(1 - \Lambda_j^+) = (1 - \Lambda_{j,0}^+)(1 - \Lambda_j^+)$ and commuting $1 - \Lambda_{j,0}^+$ until it hits Λ_{j-1}^+ , we get that this is $S_{-k}(\eta_j v)$ for any k. \Box

It follows that

Lemma B.7. For all $k \geq 1$,

$$\sum_{l=0}^{k} \|\nabla_{b}^{l}(\eta_{2}\Lambda_{1}^{+}\eta_{1}v)\| \leq C(\sum_{l=0}^{k-1} \|\nabla_{b}^{l}(\eta_{k}u)\| + \|\eta_{k+1}v\|).$$

Proof. We shall prove by induction on k that for all $j, k \ge 1$,

$$\sum_{l=0}^{k} \|\nabla_{b}^{l}(\eta_{j}\Lambda_{j-1}^{+}\eta_{j-1}v)\| \leq C(\sum_{l=0}^{k-1} \|\nabla_{b}^{l}(\eta_{j+k-2}u)\| + \|\eta_{j+k-1}v\|).$$

The case j = 2 yields the current lemma. Assume this has been proved for k - 1, and we prove the statement for k. By Lemma B.5,

$$\|\nabla_b^k(\eta_j \Lambda_{j-1}^+ \eta_{j-1} v)\| \le C(\|\nabla_b^{k-1} Z(\eta_j \Lambda_{j-1}^+ \eta_{j-1} v)\| + \|\eta_{j-1} v\|).$$

Now by Lemma B.6, one has

$$Z(\eta_j \Lambda_{j-1}^+ \eta_{j-1} v) = \eta_j \Lambda_{j-1}^+ \eta_{j-1} u + S_0(\eta_{j+1} \Lambda_j^+ \eta_j v) + S_{-(k-1)}(\eta_j v)$$

Hence one only needs to estimate $\|\nabla_b^{k-1}S_0(\eta_{j+1}\Lambda_j^+\eta_j v)\|$, which can be estimated by induction hypothesis since this involves fewer than k derivatives on $\eta_{j+1}\Lambda_j^+\eta_j v$. \Box

As a result,

Lemma B.8. For all $k \ge 0$,

$$\|\nabla_b^{k+1}v^+\| \le C(\|\nabla_b^k u^+\| + \sum_{l=0}^{k-1} \|\nabla_b^l(\eta_k u)\| + \|\eta_{k+1}v\|).$$

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Proof. By Lemma B.4,

$$\|\nabla_b^{k+1}v^+\| \le C(\|\nabla_b^k Zv^+\| + \|\eta v\|).$$

Now by Lemma B.6, one has

$$Zv^{+} = u^{+} + S_0(\eta_2 \Lambda_1^{+} \eta_1 v) + S_{-k}(\eta_1 v).$$

Hence one only needs to estimate $\|\nabla_b^k S_0(\eta_2 \Lambda_1^+ \eta_1 v)\|$, which can be estimated by Lemma B.7. \Box

Now we prove (B.8) by induction on k.

When k = 1, we need only estimate $||Zu^+||^2$. But by Lemma B.6, we have

$$u^{+} = Zv^{+} + S_{0}(\eta_{2}\Lambda_{1}^{+}\eta_{1}v) + S_{-1}(\eta_{1}v).$$

Hence

$$\begin{split} \|Zu^{+}\|^{2} \\ &= (ZZv^{+}, Zu^{+}) + O(\|Z\eta_{2}\Lambda_{1}^{+}\eta_{1}v\|\|Zu^{+}\|) + O(\|\eta_{1}v\|\|Zu^{+}\|) \\ &= (ZZv^{+}, Zu^{+}) + O(\|\eta_{1}u\|\|Zu^{+}\|) + O(\|\eta_{2}v\|\|Zu^{+}\|)) \quad \text{by Lemma B.7} \\ &= -(Zv^{+}, Z\overline{Z}u^{+}) - (Zv^{+}, iTu^{+}) + O(\|Zv^{+}\|\|\nabla_{b}u^{+}\|) \\ &+ O(\|\eta_{1}u\|\|Zu^{+}\|) + O(\|\eta_{2}v\|\|Zu^{+}\|) \\ &= -(Zv^{+}, Z\overline{Z}u^{+}) - (Zv^{+}, iTu^{+}) + O(\|\eta_{1}u\|\|\nabla_{b}u^{+}\|) \\ &+ O(\|\eta_{2}v\|\|\nabla_{b}u^{+}\|) \quad \text{by Lemma B.8} \\ &= (\overline{Z}Zv^{+}, \overline{Z}u^{+}) + (iTv^{+}, \overline{Z}u^{+}) + O(\|\nabla_{b}^{2}v^{+}\|\|u^{+}\|) + O(\|\nabla_{b}v^{+}\|(\|\nabla_{b}u^{+}\| + \|u^{+}\|)) \\ &+ O(\|\eta_{1}u\|\|\nabla_{b}u^{+}\|) + O(\|\eta_{2}v\|\|\nabla_{b}u^{+}\|) \\ &= O(\|\nabla_{b}^{2}v^{+}\|\|\overline{Z}u^{+}\|) + O(\|\nabla_{b}^{2}v^{+}\|\|u^{+}\|) + O(\|\nabla_{b}v^{+}\|(\|\nabla_{b}u^{+}\| + \|u^{+}\|)) \\ &+ O(\|\eta_{1}u\|\|\nabla_{b}u^{+}\|) + O(\|\eta_{2}v\|\|\nabla_{b}u^{+}\|) \\ &\leq C(\varepsilon\|\nabla_{b}^{2}v^{+}\|^{2} + \varepsilon^{-1}\|\overline{Z}u^{+}\|^{2} + \varepsilon^{-1}\|u^{+}\|^{2} + \varepsilon^{-1}\|\nabla_{b}v^{+}\|^{2} \\ &+ \varepsilon^{-1}\|\eta_{1}u\|^{2} + \varepsilon^{-1}\|\eta_{2}v\|^{2}). \end{split}$$

Absorbing $C\varepsilon ||Zu^+||^2$ to the left hand side, we get

$$||Zu^{+}||^{2} \leq C(\varepsilon ||\nabla_{b}^{2}v^{+}||^{2} + \varepsilon^{-1} ||\overline{Z}u^{+}||^{2} + \varepsilon^{-1} ||\nabla_{b}v^{+}||^{2} + \varepsilon^{-1} ||\eta_{1}u||^{2} + \varepsilon^{-1} ||\eta_{2}v||^{2}).$$

Now by Lemma B.8, one estimates $\|\nabla_b^2 v^+\|$ and $\|\nabla_b v^+\|$:

$$\|\nabla_b^2 v^+\| \le C \left(\|\nabla_b u^+\| + \|\eta_1 u\| + \|\eta_2 v\| \right),$$

$$\|\nabla_b v^+\| \le C (\|u^+\| + \|\eta_1 v\|) \le C (\|\eta_1 u\| + \|\eta_2 v\|).$$

Together, we get

$$||Zu^+||^2 \le C\left(||\overline{Z}u^+||^2 + ||\eta_1 u|| + ||\eta_2 v||\right)$$

as desired.

Next, to prove (B.8) for a general k, we prove the following statement by induction on m for all $0 \le m \le \lfloor \frac{k}{2} \rfloor$ and $\varepsilon > 0$:

$$\|T^{m}\nabla_{b}^{k-2m}u^{+}\| \leq \varepsilon \|\nabla_{b}^{k}u^{+}\| + C_{\varepsilon}\left(\|\nabla_{b}^{k-1}\overline{Z}u^{+}\| + \sum_{l=0}^{k-1} \|\nabla_{b}^{l}(\eta_{k}u)\| + \|\eta_{k+1}v\|\right)$$
(B.10)

In fact the case m = 0 readily implies (B.8) for k.

Again we begin from the biggest possible value of m. Suppose k - 2m = 0. Then we need to estimate $||T^m u^+||$. Now

$$\begin{split} \|T^{m}u^{+}\|^{2} &= (T^{m}Zv^{+}, T^{m}u^{+}) + O((\|\nabla_{b}^{k}(\eta_{2}\Lambda_{1}^{+}\eta_{1}v)\| + \|\eta_{1}v\|)\|T^{m}u^{+}\|) \\ &= -(T^{m}v^{+}, T^{m}\overline{Z}u^{+}) + O(\sum_{l=0}^{k}\|\nabla_{b}^{l}v^{+}\|\|T^{m}u^{+}\|) + O(\|T^{m}v^{+}\|\sum_{l=0}^{k}\|\nabla_{b}^{l}u^{+}\|) \\ &+ O((\sum_{l=0}^{k-1}\|\nabla_{b}^{l}(\eta_{k}u)\| + \|\eta_{k+1}v\|)\|T^{m}u^{+}\|) \\ &= (\nabla_{b}T^{m}v^{+}, \nabla_{b}T^{m-1}\overline{Z}u^{+}) + O(\sum_{l=0}^{k}\|\nabla_{b}^{l}v^{+}\|\|T^{m}u^{+}\|) + O(\|T^{m}v^{+}\|\sum_{l=0}^{k}\|\nabla_{b}^{l}u^{+}\|) \\ &+ O((\sum_{l=0}^{k-1}\|\nabla_{b}^{l}(\eta_{k}u)\| + \|\eta_{k+1}v\|)\|T^{m}u^{+}\|) \end{split}$$

the last line following by writing one of the T's in $T^m \overline{Z} u^+$ as a commutator and integrating by parts. Hence

$$||T^{m}u^{+}||^{2} \leq \varepsilon ||\nabla_{b}^{k+1}v^{+}||^{2} + \varepsilon ||\nabla_{b}^{k}u^{+}||^{2} + C_{\varepsilon}(||\nabla_{b}^{k-1}\overline{Z}u^{+}||^{2} + \sum_{l=0}^{k} ||\nabla_{b}^{l}v^{+}||^{2} + \sum_{l=0}^{k-1} ||\nabla_{b}^{l}(\eta_{k}u)||^{2} + ||\eta_{k+1}v||^{2}).$$

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Now we invoke Lemma B.8 to estimate $\|\nabla_b^{k+1}v^+\|$ and $\sum_{l=0}^k \|\nabla_b^l v^+\|$. Together,

$$||T^{m}u^{+}||^{2} \leq \varepsilon ||\nabla_{b}^{k}u^{+}||^{2} + C_{\varepsilon} \left(||\nabla_{b}^{k-1}\overline{Z}u^{+}||^{2} + \sum_{l=0}^{k-1} ||\nabla_{b}^{l}(\eta_{k}u)||^{2} + ||\eta_{k+1}v||^{2} \right)$$

which implies the desired estimate for $T^m u^+$.

Next, we estimate $||T^m Z^{k-2m} u^+||$ for any $0 \le m \le \lfloor \frac{k}{2} \rfloor$, if k - 2m > 0. Then

$$\begin{split} \|T^{m}Z^{k-2m}u^{+}\|^{2} &= (T^{m}Z^{k-2m+1}v^{+}, T^{m}Z^{k-2m}u^{+}) + O((\|\nabla_{b}^{k}(\eta_{2}\Lambda_{1}^{+}\eta_{1}v)\| + \|\eta_{1}v\|)\|T^{m}Z^{k-2m}u^{+}\|) \\ &= -(T^{m}Z^{k-2m}v^{+}, T^{m}Z^{k-2m}\overline{Z}u^{+}) - (k-2m)(T^{m}Z^{k-2m}v^{+}, iT^{m+1}Z^{k-2m-1}u^{+}) \\ &+ O(\sum_{l=0}^{k}\|\nabla_{b}^{l}v^{+}\|\sum_{j=0}^{k}\|\nabla_{b}^{j}u^{+}\|) + O((\sum_{l=0}^{k-1}\|\nabla_{b}^{l}(\eta_{k}u)\| + \|\eta_{k+1}v\|)\|T^{m}Z^{k-2m}u^{+}\|) \end{split}$$

Now in the first term, we split off one Z in $T^m Z^{k-2m} \overline{Z}$ and integrate by parts. Also, in the second term, we split off one T in $T^{m+1}Z^{k-2m-1}$ and integrate by parts; then we split off one Z in $T^m Z^{k-2m}$ and integrate by parts. We get

$$\begin{split} \|T^{m}Z^{k-2m}u^{+}\|^{2} &= (\overline{Z}T^{m}Z^{k-2m}v^{+}, T^{m}Z^{k-2m-1}\overline{Z}u^{+}) + (k-2m)(iT^{m+1}Z^{k-2m-1}v^{+}, T^{m}Z^{k-2m}u^{+}) \\ &+ O(\|\nabla_{b}^{k+1}v^{+}\|\sum_{l=0}^{k-1}\|\nabla_{b}^{l}u^{+}\|) + O(\sum_{l=0}^{k}\|\nabla_{b}^{l}v^{+}\|\sum_{j=0}^{k}\|\nabla_{b}^{l}u^{+}\|) \\ &+ O((\sum_{l=0}^{k-1}\|\nabla_{b}^{l}(\eta_{k}u)\| + \|\eta_{k+1}v\|)\|T^{m}Z^{k-2m}u^{+}\|) \\ &\leq O(\|\nabla_{b}^{k+1}v^{+}\|\|\nabla_{b}^{k-1}\overline{Z}u^{+}\|) + O(\|\nabla_{b}^{k+1}v^{+}\|\|\sum_{l=0}^{k-1}\|\nabla_{b}^{l}u^{+}\|) \\ &+ O(\sum_{l=0}^{k}\|\nabla_{b}^{l}v^{+}\|\sum_{j=0}^{k}\|\nabla_{b}^{l}u^{+}\|) + O((\sum_{l=0}^{k-1}\|\nabla_{b}^{l}(\eta_{k}u)\| + \|\eta_{k+1}v\|)\|T^{m}Z^{k-2m}u^{+}\|) \end{split}$$

As a result,

$$\begin{aligned} \|T^{m}Z^{k-2m}u^{+}\|^{2} \\ &\leq \varepsilon \|\nabla_{b}^{k+1}v^{+}\|^{2} + \varepsilon \|\nabla_{b}^{k}u^{+}\|^{2} \\ &+ C_{\varepsilon} \left(\|\nabla_{b}^{k-1}\overline{Z}u^{+}\|^{2} + \sum_{l=0}^{k} \|\nabla_{b}^{l}v^{+}\|^{2} + \sum_{l=0}^{k-1} \|\nabla_{b}^{l}(\eta_{k}u)\|^{2} + \|\eta_{k+1}v\|^{2} \right) \end{aligned}$$

and the desired estimate follows upon invoking Lemma B.8.

Now suppose we have proved (B.10) for all strictly bigger m's, and want to prove the inequality for m. Then we need to estimate $||T^m \nabla_b^{k-2m} u^+||$. If one of the ∇_b is \overline{Z} , we commute until that \overline{Z} hits u^+ , obtaining an error that has more T in it, which one can estimate by the induction hypothesis. Otherwise all ∇_b are Z's, and the estimate follows from what we have proved above. This completes the proof of (B.8).

Now putting (B.1), (B.6) and (B.8) together, and remembering that $\eta u = u^- + u^0 + u^+$, we get

$$\|\nabla_{b}^{k}(\eta u)\| \leq C_{k} \left(\|\nabla_{b}^{k-1}\overline{Z}(\eta u)\| + \sum_{l=0}^{k-1} \|\nabla_{b}^{l}(\eta_{k}u)\| + \|\eta_{k+1}v\| \right).$$

But $\sum_{l=0}^{k-1} \|\nabla_b^l(\eta_k u)\|$ involves fewer than k derivatives of $\eta_k u$, and can be estimated if we iterate the above inequality. In fact

$$\sum_{l=0}^{k-1} \|\nabla_b^l(\eta_k u)\| \le C \left(\sum_{l=0}^{k-2} \|\nabla_b^l \overline{Z}(\tilde{\eta} u)\| + \|\tilde{\eta} u\| + \|\tilde{\eta} v\| \right).$$

It follows that

$$\|\nabla_b^k(\eta u)\| \le C_k \left(\sum_{l=0}^{k-1} \|\nabla_b^l \overline{Z}(\tilde{\eta} u)\| + \|\tilde{\eta} v\|\right).$$

Using the interpolation inequality in Lemma B.3, one obtains the desired inequality in Proposition B.1.

Finally, we come back to the remark we made after the statement of Proposition B.1. We remark that if we have $Zv + \alpha v = u$ instead of Zv = u, where α is a fixed C^{∞} function on Q_2 , then the above theorem still holds. See the end of this section for a discussion about that. The key there is to observe that Lemma B.6 above holds under this modified assumption as well. In fact, then

$$\begin{split} \eta_{j}\Lambda_{j-1}^{+}\eta_{j-1}u &= \eta_{j}\Lambda_{j-1}^{+}\eta_{j-1}(Z+\alpha)v\\ &= \eta_{j}\Lambda_{j-1}^{+}\eta_{j-1}(Z+\alpha)(\eta_{j}v)\\ &= \eta_{j}\Lambda_{j-1}^{+}\eta_{j-1}(Z+\alpha)(\eta_{j+1}\Lambda_{j}^{+}\eta_{j}v)\\ &\quad + \eta_{j}\Lambda_{j-1}^{+}\eta_{j-1}(Z+\alpha)(\eta_{j+1}(1-\Lambda_{j}^{+})\eta_{j}v). \end{split}$$

The α in the first term contributes only $S_0(\eta_{j+1}\Lambda_j^+\eta_j v)$, while the α in the last term contributes only $S_{-k-1}(\eta_j v)$.

References

R. Beals, P. Greiner, Calculus on Heisenberg Manifolds, Ann. of Math. Stud., vol. 119, Princeton University Press, Princeton, NJ, 1988, MR953082 (89m:35223).

- H.P. Boas, M.-C. Shaw, Sobolev estimates for the Lewy operator on weakly pseudoconvex boundaries, Math. Ann. 274 (2) (1986) 221–231, http://dx.doi.org/10.1007/BF01457071, MR838466 (87i:32029).
- [3] L. Boutet de Monvel, J. Sjöstrand, Sur la singularité des noyaux de Bergman et de Szegő, in: Journées: Équations aux Dérivées Partielles de Rennes, 1975, Soc. Math. France, Paris, Astérisque 34–35 (1976) 123–164 (in French), MR0590106 (58 #28684).
- [4] S. Chanillo, H.-L. Chiu, P. Yang, Embeddability for 3-dimensional Cauchy–Riemann manifolds and CR Yamabe invariants, Duke Math. J. 161 (15) (2012) 2909–2921, http://dx.doi.org/10.1215/ 00127094-1902154, MR2999315.
- [5] J.-H. Cheng, A. Malchiodi, P. Yang, A positive mass theorem in three dimensional Cauchy–Riemann geometry, preprint.
- [6] M. Christ, Regularity properties of the $\bar{\partial}_b$ equation on weakly pseudoconvex CR manifolds of dimension 3, J. Amer. Math. Soc. 1 (3) (1988) 587–646, http://dx.doi.org/10.2307/1990950, MR928903 (89e:32027).
- M. Christ, Pointwise estimates for the relative fundamental solution of ∂
 _b, Proc. Amer. Math. Soc. 104 (3) (1988) 787–792, http://dx.doi.org/10.2307/2046793, MR929407 (89d:35126).
- [8] C. Fefferman, The Bergman kernel and biholomorphic mappings of pseudoconvex domains, Invent. Math. 26 (1974) 1–65, MR0350069 (50 #2562).
- C.L. Fefferman, J.J. Kohn, Estimates of kernels on three-dimensional CR manifolds, Rev. Mat. Iberoam. 4 (3–4) (1988) 355–405, http://dx.doi.org/10.4171/RMI/78, MR1048582 (91h:32013).
- [10] G.B. Folland, E.M. Stein, Estimates for the $\bar{\partial}_b$ complex and analysis on the Heisenberg group, Comm. Pure Appl. Math. 27 (1974) 429–522, MR0367477 (51 #3719).
- [11] P.C. Greiner, E.M. Stein, Estimates for the ∂-Neumann Problem, Math. Notes, vol. 19, Princeton University Press, Princeton, NJ, 1977, MR0499319 (58 #17218).
- [12] C.-Y. Hsiao, Projections in several complex variables, Mém. Soc. Math. Fr. (N.S.) 123 (2010) 131 (in English, with English and French summaries), MR2780123 (2011m:32004).
- [13] C.-Y. Hsiao, P.-L. Yung, The tangential Cauchy–Riemann complex on the Heisenberg group via Conformal Invariance, Bull. Inst. Math. Acad. Sin. (N.S.) 8 (3) (2013) 359–375.
- [14] F. Jean, Control of nonholonomic systems and sub-Riemannian geometry, preprint.
- [15] K.D. Koenig, On maximal Sobolev and Hölder estimates for the tangential Cauchy–Riemann operator and boundary Laplacian, Amer. J. Math. 124 (1) (2002) 129–197, MR1879002 (2002m:32061).
- [16] J.J. Kohn, Estimates for $\bar{\partial}_b$ on pseudoconvex CR manifolds, in: Pseudodifferential Operators and Applications, Notre Dame, Ind., 1984, in: Proc. Sympos. Pure Math., vol. 43, Amer. Math. Soc., Providence, RI, 1985, pp. 207–217, MR812292 (87c:32025).
- [17] J.J. Kohn, The range of the tangential Cauchy–Riemann operator, Duke Math. J. 53 (2) (1986) 525–545, http://dx.doi.org/10.1215/S0012-7094-86-05330-5, MR850548 (87m:32041).
- [18] J.J. Kohn, H. Rossi, On the extension of holomorphic functions from the boundary of a complex manifold, Ann. of Math. (2) 81 (1965) 451–472, MR0177135 (31 #1399).
- [19] X. Ma, G. Marinescu, Holomorphic Morse Inequalities and Bergman Kernels, Progr. Math., vol. 254, Birkhäuser Verlag, Basel, 2007, MR2339952 (2008g:32030).
- [20] M. Machedon, Estimates for the parametrix of the Kohn Laplacian on certain domains, Invent. Math. 91 (2) (1988) 339–364, http://dx.doi.org/10.1007/BF01389371, MR922804 (89d:58118).
- [21] M. Machedon, Szegő kernels on pseudoconvex domains with one degenerate eigenvalue, Ann. of Math. (2) 128 (3) (1988) 619–640, http://dx.doi.org/10.2307/1971438, MR970613 (89i:32043).
- [22] A. Nagel, E.M. Stein, Lectures on Pseudodifferential Operators: Regularity Theorems and Applications to Noelliptic Problems, Math. Notes, vol. 24, Princeton University Press, Princeton, NJ, 1979, MR549321 (82f:47059).
- [23] A. Nagel, E.M. Stein, S. Wainger, Balls and metrics defined by vector fields. I. Basic properties, Acta Math. 155 (1–2) (1985) 103–147, http://dx.doi.org/10.1007/BF02392539, MR793239 (86k:46049).
- [24] A. Nagel, J.-P. Rosay, E.M. Stein, S. Wainger, Estimates for the Bergman and Szegő kernels in certain weakly pseudoconvex domains, Bull. Amer. Math. Soc. (N.S.) 18 (1) (1988) 55–59, http://dx.doi.org/10.1090/S0273-0979-1988-15598-X, MR919661 (89a:32025).
- [25] A. Nagel, J.-P. Rosay, E.M. Stein, S. Wainger, Estimates for the Bergman and Szegő kernels in C², Ann. of Math. (2) 129 (1) (1989) 113–149, http://dx.doi.org/10.2307/1971487, MR979602 (90g:32028).
- [26] L.P. Rothschild, E.M. Stein, Hypoelliptic differential operators and nilpotent groups, Acta Math. 137 (3-4) (1976) 247-320, MR0436223 (55 #9171).

[27] E.M. Stein, Harmonic Analysis: Real-variable Methods, Orthogonality, and Oscillatory Integrals, Princeton Math. Ser., vol. 43, Princeton University Press, Princeton, NJ, 1993, with the assistance of Timothy S. Murphy; Monographs in Harmonic Analysis, III, MR1232192 (95c:42002).