

EXTENSION OF CR STRUCTURES ON PSEUDOCONVEX CR MANIFOLDS WITH ONE DEGENERATE EIGENVALUE

SANGHYUN CHO

ABSTRACT. Let \overline{M} be a smoothly bounded orientable pseudoconvex CR manifold of finite type with at most one degenerate eigenvalue. Then we extend the given CR structure on M to an integrable almost complex structure on the concave side of M . Therefore we may regard M as the boundary of a complex manifold.

1. INTRODUCTION

Suppose that \overline{M} is an abstract smoothly bounded orientable CR manifold of dimension $2n - 1$ with a given integrable CR structure \mathcal{S} of dimension $n - 1$. Since \overline{M} is orientable, there are a smooth real nonvanishing 1-form η and a smooth real vector field X_0 on \overline{M} so that $\eta(X) = 0$ for all $X \in \mathcal{S}$ and $\eta(X_0) = 1$. We define the Levi form of \mathcal{S} on \overline{M} by $i\eta([X', \overline{X}''])$, $X', X'' \in \mathcal{S}$. We may assume that $\overline{M} \subset \widetilde{M}$, in C^∞ sense, where \widetilde{M} is a smooth manifold.

In [5], Catlin has studied an extension problem of a given CR structure on M to an integrable almost complex structure on a $2n$ -dimensional manifold Ω with boundary, $b\Omega$, so that the extension is smooth up to the boundary and so M lies in $b\Omega$. Under certain conditions on the Levi form (cf. [5, Theorem 1.1, Theorem 1.3]), this shows that an abstract CR manifold can be locally embedded in \mathbb{C}^n [1,13,16].

In this paper, we study an extension problem of a given CR structure on M when M is a pseudoconvex CR manifold of finite type with one degenerate eigenvalue and $\dim_{\mathbb{R}} M = 2n - 1$. For a given positive continuous function g on M , where $g = 0$ on bM , the boundary of M , we define

$$S_g^+ = \{(x, t) \in M \times [0, \infty) ; 0 \leq t \leq g(x)\}.$$

Then our main result is the following theorem:

Theorem 1.1. *Let $(\overline{M}, \mathcal{S})$ be a smoothly bounded pseudoconvex CR manifold of finite type with at most one degenerate eigenvalue and $\dim_{\mathbb{R}} M = 2n - 1$. Then there exists a positive continuous function g on M and a smooth integrable almost complex structure \mathcal{L} on S_g^+ such that for all $x \in M$, $\mathcal{L}_{(x,0)} \cap \mathbb{C}TM = \mathcal{S}_x$. Furthermore, if $\mathcal{J}_{\mathcal{L}} : TS_g^+ \rightarrow TS_g^+$ is the map associated with the complex structure \mathcal{L} , then $dt(\mathcal{J}_{\mathcal{L}}(X_0)) < 0$ at all points of $M_0 = \{(x, 0) ; x \in M\}$.*

Note that we extend the given CR structure on M to the concave side (instead of convex side) of M . When $\dim_{\mathbb{R}} M = 3$, the author proved the same result when \overline{M} is a pseudoconvex

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CR manifold of finite type [8]. We also note that if M is strongly pseudoconvex, this case was handled in [5, Theorem 1.1].

In general, Theorem 1.1 does not imply that the given CR structure can be locally embedded in \mathbb{C}^n [12]. On the other hand, a theorem of Newlander-Nirenberg [14] shows that an integrable almost complex manifold is a complex manifold. Therefore if we let (M, \mathcal{S}) be as in Theorem 1.1, then we have the following corollaries.

Corollary 1.2. *We may regard (M, \mathcal{S}) as a boundary of a complex manifold.*

Corollary 1.3. *Let D be a complex manifold with C^∞ boundary, bD , and $\dim_{\mathbb{C}} D = n$. Suppose that the almost complex structure on D extends smoothly to a manifold $\overline{M} \subset bD$, where \overline{M} is a smoothly bounded pseudoconvex CR manifold of finite type and the Levi-form of \overline{M} has one degenerate eigenvalue and $\dim_{\mathbb{R}} M = 2n - 1$. Then D can be embedded into a larger complex manifold Ω so that M lies in the interior of Ω as a real hypersurface.*

Remark 1.4. In [6], the author showed that any smooth compact pseudoconvex complex manifold \overline{D} of finite type with $\dim_{\mathbb{C}} D = n$, $n \geq 2$, can be embedded into a larger complex manifold Ω . Corollary 1.3 is a generalization of this result to some special non-compact complex manifolds.

Remark 1.5. If (M, \mathcal{S}) has at least three positive eigenvalues, Catlin [5] has extended the given CR-structure of \mathcal{S} to the pseudoconvex side of \overline{M} [5, Theorem 1.1]. If M is also pseudoconvex, this result implies the local embedding of CR-structures in \mathbb{C}^n .

In [5], Catlin has introduced certain nonlinear equations which stem from deformation theory of an almost complex structure (Section 2). The linearized forms of these equations are simply the $\overline{\partial}$ -operator from $\Lambda^{0,1} \otimes T^{1,0}$ to $\Lambda^{0,2} \otimes T^{1,0}$ (Section 2). To overcome difficulties in subelliptic estimates for $\overline{\partial}$ near bM , we choose a Hermitian metric on S_g^+ so that S_g^+ takes on the form $S_\varepsilon = M \times [0, \varepsilon]$. To this end, we choose, for each $x_0 \in M$, a noneuclidean ball that is of size $\delta = g(x_0)$ in the transverse holomorphic direction and of size $\delta^{1/2}$ in strongly pseudoconvex tangential holomorphic directions, and of size $\tau(x_0, \delta)$ in the weakly pseudoconvex tangential holomorphic direction. We choose the metric so that the unit ball about $x_0 \in M$ corresponds to the above noneuclidean ball with $\delta = g(x_0)$.

To show that $\tau(x_0, \delta)$ is invariantly defined (i.e., independent of coordinate functions), we choose special coordinates defined near $x_0 \in M$ (Proposition 3.1). These change of coordinates shall have an independent interest in studying weakly pseudoconvex CR manifolds of finite type. In terms of these special coordinates, the weakly pseudoconvex tangential holomorphic vector field L_1 has a special representation so that we can define another quantity $\mu(x, \delta)$, which is a smooth function of δ and x and it is obtained by taking successive brackets of L_1 and \overline{L}_1 and hence defined invariantly. A technical difficulty is to show that the brackets mixed with L_1 and the strongly pseudoconvex tangential vector fields are not major terms in determining $\mu(x, \delta)$. Then we show that $\tau(x, \delta) \approx \mu(x, \delta)$ (Proposition 3.5), and hence $\tau(x, \delta)$ is defined invariantly. To get an $0 < \varepsilon \leq 1/2$ subelliptic estimates (other than $1/2$ estimates of Catlin [5]), we also need some precise estimates for $\eta([L_1, \overline{L}_k])$, $2 \leq k \leq n$. This difficulty comes from the fact that we have a Dirichlet condition on M_0 and hence we need to control the boundary integral terms, which were not occurred in dimension three case, on M_0 (In usual $\overline{\partial}$ -estimates we deal with forms which vanish on M_0 and hence there are no boundary integral terms). Section 3, 4 and 5 contain these estimates in detail.

After this, we technically construct a family of plurisubharmonic functions with large Hessian using the properties of $\tau(x, \delta)$ and $\mu(x, \delta)$. In performing the subelliptic estimates, we use $(n - 2)$ -positive eigenvalue conditions on M_σ to handle the boundary integral terms on M_σ ,

and we use the existence of a family of plurisubharmonic functions with large Hessian to handle the components vanishing on M_σ . This will give us uniform $1/m$ subelliptic estimates for $\bar{\partial}$ on each non-euclidean ball. Then we get the estimates, so-called ‘‘tame estimates’’, which are required in the Nash-Moser theorem [14] for the approximate solution to the linearized equation.

2. DEFORMATION OF ALMOST COMPLEX STRUCTURES

Let (M, \mathcal{S}) be a CR manifold as in Section 1 and set $\Omega = M \times (-1, 1)$. In this section we extend the given CR structure \mathcal{S} on M to an almost complex manifold (Ω, \mathcal{L}) , and study a deformation problem of the almost complex structure \mathcal{L} on Ω so that the new (deformed) almost complex structure is integrable (or close to be integrable).

Assume that \mathcal{L} is an almost complex structure on Ω . Let A be a smooth section of $\Gamma^1(\mathcal{L}) = \Lambda^{0,1}(\mathcal{L}) \otimes \mathcal{L}$, where $\Lambda^{0,1}(\mathcal{L})$ denotes the set of $(0, 1)$ forms with respect to \mathcal{L} . Observe that if A is sufficiently small, then the bundle $\mathcal{L}^A = \{L + \bar{A}(L) ; L \in \mathcal{L}\}$ defines a new almost complex structure and if \bar{L}' and \bar{L}'' are sections of $\bar{\mathcal{L}}$, then $\bar{L}' + A(\bar{L}')$ and $\bar{L}'' + A(\bar{L}'')$ are sections of $\bar{\mathcal{L}}^A$. Similarly, if ω is a section of $\Lambda^{1,0}(\mathcal{L})$, then $\omega - A^*\omega$ is a section of $\Lambda^{1,0}(\mathcal{L}^A)$, where the adjoint A^* maps from $\Lambda^{1,0}(\mathcal{L})$ to $\Lambda^{0,1}(\mathcal{L})$ and is defined by

$$(2.1) \quad (A^*\omega)(\bar{L}) = \omega(A(\bar{L}))$$

for all $\bar{L} \in \bar{\mathcal{L}}$ and $\omega \in \Lambda^{1,0}$. We want to choose A so that

$$(\omega - A^*\omega)([L' + A(L'), L'' + A(L'')]) = 0.$$

By linearizing, i.e., by ignoring terms where A or A^* appear more than once, we obtain

$$(2.2) \quad \omega([L', A(L'')]) + \omega([A(L'), L'']) - A^*\omega([L', L'']) = -\omega([L', L'']).$$

Let $L = L' + L''$ denote the decomposition of a vector $L \in \mathbb{C}T_z$ where $L' \in \mathcal{L}_z$ and $L'' \in \bar{\mathcal{L}}_z$.

For sections \bar{L}_1, \bar{L}_2 of $\bar{\mathcal{L}}$, we define

$$(2.3) \quad (D_2A)(\bar{L}_1, \bar{L}_2) = [\bar{L}_1, A(\bar{L}_2)]' - [\bar{L}_2, A(\bar{L}_1)]' - A([\bar{L}_1, \bar{L}_2]'').$$

Note that this definition is linear in \bar{L}_1 and \bar{L}_2 , and hence D_2A is a section of $\Gamma^2 = \Lambda^{0,2}(\mathcal{L}) \otimes \mathcal{L}$. It follows from (2.1) and (2.3) that (2.2) is equivalent to the equation

$$(2.4) \quad D_2A = -F,$$

where F is a section of Γ^2 defined by

$$(2.5) \quad F(\bar{L}_1, \bar{L}_2) = [\bar{L}_1, \bar{L}_2]'$$

Note that F measures to what extent \mathcal{L} fails to be integrable. If \mathcal{L} defines a CR structure on M and if we want \mathcal{L}_A to define the same CR structure on M , then this means that A must satisfy $A(\bar{L}') = 0$ on M whenever \bar{L}' is a section of $\bar{\mathcal{L}}$ that is tangent to M . This is a Dirichlet condition on some of the components of the solution (2.4).

If we define $D_3 : \Gamma^2 \rightarrow \Gamma^3 = \Lambda^{0,3}(\mathcal{L}) \otimes \mathcal{L}$ by

$$(2.6) \quad \begin{aligned} D_3B(\bar{L}_1, \bar{L}_2, \bar{L}_3) = & [\bar{L}_1, B(\bar{L}_2, \bar{L}_3)]' - [\bar{L}_2, B(\bar{L}_1, \bar{L}_3)]' + [\bar{L}_3, B(\bar{L}_1, \bar{L}_2)]' \\ & - B([\bar{L}_1, \bar{L}_2]'', \bar{L}_3) + B([\bar{L}_1, \bar{L}_3]'', \bar{L}_2) - B([\bar{L}_2, \bar{L}_3]'', \bar{L}_1) \end{aligned}$$

for $B \in \Gamma^2$, then it follows that $D_3F = 0$ [5, Lemma 3.2]. Then we have the following formal solution of the extension problem [5, Theorem 4.1].

Theorem 2.1. *Suppose that M is an orientable CR manifold of dimension $2n - 1$ such that the CR dimension equals $n - 1$. Then there exists an almost complex structure \mathcal{L}^* on $\Omega = M \times (-1, 1)$ such that \mathcal{L}^* is an extension of the CR structure on M , and such that it is integrable to infinite order at M in the sense that if ω is a section of $\Lambda^{1,0}(\mathcal{L}^*)$ and \bar{L}_1, \bar{L}_2 are sections of $\bar{\mathcal{L}}^*$, then $\omega([\bar{L}_1, \bar{L}_2])$ vanishes to infinite order along M .*

Let M and Ω be as in Theorem 2.1. The next theorem shows that the above formal extension is essentially unique [5, Theorem 4.2].

Theorem 2.2. *Suppose that \mathcal{L} and \mathcal{X} are almost complex structures on Ω that extend the CR-structure on $M_0 = \{(x, 0) ; x \in M\}$, and that are integrable to infinite order on M_0 . Then there exists a diffeomorphism G of Ω onto itself that is the identity when $t = 0$ and such that $G_*\mathcal{X}$ approximates \mathcal{L} to infinite order near M_0 in the sense that if X is a section of \mathcal{L} , then G_*X differs from a section of \mathcal{L} by a vector field which vanishes to infinite order on M_0 .*

By Theorem 2.1, we have an almost complex structure \mathcal{L}^* on Ω , that is integrable to infinite order along $M_0 = \{(x, 0) ; x \in M\}$. Let η be a smooth non-vanishing one form on M that satisfies $\eta(L) = 0$ for all $L \in \mathcal{S}_x$, $x \in M$, and that defines the Levi form of M as in Section 1. We can clearly extend η to all of Ω so that it still annihilates $\mathcal{S}_{(x,t)}$ for all $(x, t) \in \Omega$, where $\mathcal{S}_{(x,t)}$ now denotes the space of vectors in $\mathcal{L}^*_{(x,t)}$ that are tangent to the level set of the auxiliary coordinate t . Then we have the following theorem which is a formal solution of local embedding problem. One can refer a proof to, for example, [3, Proposition 3].

Theorem 2.3. *Let $x_0 \in M$. Then there are a small neighborhood U of x_0 and a constant $c > 0$ so that for each $x \in M \cap U$, there are almost holomorphic functions f_1, \dots, f_n defined on \bar{U} so that if $F_x = (f_1, \dots, f_n)$, then $F_x(x) = 0$ and*

- (a) $|dF_x| \geq c$ on \bar{U} , and
- (b) $\bar{L}f_j$ vanishes to infinite order at x_0 for each $L \in \mathcal{L}$.

Remark 2.4. Suppose $L \in \mathcal{S}_{(x,0)}$. Then $(F_{x_0})_*L$ differs from a section of $T^{1,0}(F_{x_0}(M))$ by a vector field which vanishes to infinite order at 0. Therefore the image $F_{x_0}(M)$ is a smooth real hypersurface in \mathbb{C}^n with defining function given by $r(w) = t \circ F_{x_0}^{-1}(w)$.

In order to define the type of $x_0 \in M$, we use the (almost) holomorphic function F_x constructed in Theorem 2.3:

Definition 2.5. Let (x_0, U, F_{x_0}) be as in Theorem 2.3. Then we define the type of x_0 is equal to the type of $F_{x_0}(x_0) = 0 \in \mathbb{C}^n$ in the sense of D'Angelo [9].

Set $T(x_0) =$ the type of $x_0 \in M$, and set

$$T(\bar{M}) = \max\{T(x_0) ; x_0 \in \bar{M}\} = m.$$

Under the assumption that the Levi-form of M has $(n - 2)$ -positive eigenvalues, we may assume that m is an even integer.

Let us take (Ω, \mathcal{L}^*) constructed in Theorem 2.1. Choose a smooth real vector field X_0 on Ω that satisfies $X_0t \equiv 0$ and $\eta(X_0) \equiv 1$ in Ω . Set $Y_0 = -\mathcal{J}_{\mathcal{L}^*}(X_0)$ so that $X_0 + iY_0$ is a section of \mathcal{L}^* that is transverse to the level set of t . Let $G : \Omega \rightarrow \Omega$ be a diffeomorphism such that G fixes M_0 and

$$G_*Y_0|_{(x,0)} = \frac{\partial}{\partial t}|_{(x,0)}, \quad x \in M.$$

Since M is orientable, we may assume that $dt(\mathcal{J}_{\mathcal{L}^*}(X_0)) < 0$. Thus $dt(Y_0) > 0$ along M_0 , which shows that G preserves the sides of M_0 , i.e., G maps $\Omega^+ = \{(x, t) ; 0 \leq t < 1\}$ into itself. If we set $\mathcal{L}^0 = G_*\mathcal{L}^*$, then clearly $\tilde{Z} = -iG_*(X_0 + iY_0)$ is a section of \mathcal{L}^0 such that along M_0 ,

$$\tilde{Z} = -iX_0 + \frac{\partial}{\partial t}.$$

We write $\tilde{Z} = \tilde{X} + g(x, t)\partial/\partial t$ where $\tilde{X}t \equiv 0$, and set $L_n = g^{-1}\tilde{Z}$. Then $L_n = \partial/\partial t + X$, where $Xt \equiv 0$. We fix a smooth metric $\langle \cdot, \cdot \rangle_0$ that is Hermitian with respect to the structure \mathcal{L}^0 on Ω . Note that along M , we have $L_n = \partial/\partial t - iX_0$ and $dt = (1/2)(dt + i\eta) + (1/2)(dt - i\eta)$, which implies that $\partial t = (1/2)(dt + i\eta)$. Hence $\partial t(L) = (1/2)dt(L) + (i/2)\eta(L)$ and

$$\partial t([X_1, \bar{X}_2]) = \frac{i}{2}\eta([X_1, \bar{X}_2])$$

for all $X_1, X_2 \in \mathcal{S}_{(x,t)}$, along M .

3. CONSTRUCTION OF PLURISUBHARMONIC FUNCTIONS

Let M, Ω, X_0 and \mathcal{L}^0 be as in Section 2. In this section, we will construct a family of plurisubharmonic functions with large Hessian (Theorem 3.6). For this purpose, we first construct special coordinate functions defined in a neighborhood of $x_0 \in M$ so that the coefficients of the weakly pseudoconvex tangential holomorphic vector fields satisfy some necessary estimates in new coordinates (Lemma 3.7 and Proposition 3.8).

Assume that $\tilde{x}_0 \in M$. Then there are coordinate functions x_1, \dots, x_{2n} defined on a neighborhood \bar{U} of \tilde{x}_0 with the property that $x_{2n} = t$ and that

$$x_k(x', t) = x_k(x', 0), \quad k < 2n,$$

for $(x', t) \in \bar{U}$, and that

$$\frac{\partial}{\partial x_{2n-1}} = -X_0$$

at all points of $\bar{U} \cap M$. We take an orthonormal frame $\{\tilde{L}_1, \dots, \tilde{L}_n\}$ of \mathcal{L}^0 defined on \bar{U} . Let $x_0 \in M \cap U$ be fixed for a moment. If \tilde{L}_j is replaced by $L_j = \sum_{k=1}^{n-1} U_{jk}\tilde{L}_k$ where $U = (U_{jk})$ is a suitably chosen unitary matrix so that

$$(3.1) \quad \frac{i}{2}\eta([L_j, \bar{L}_k])(x_0) = \delta_{jk}d_j(x_0) := d_{j,k}(x_0), \quad 1 \leq j, k \leq n-1,$$

where $0 \leq d_1 \leq d_2 \leq \dots \leq d_{n-1}$, and $d_j(x_0)$ is a smooth function defined on \bar{U} satisfying $d_2(x) \geq d_0 > 0$ on \bar{U} for a uniform constant $d_0 > 0$, while δ_{jk} is a Dirac delta function.

There is an affine transformation $C_{x_0} : \mathbb{R}^{2n} \rightarrow \mathbb{R}^{2n}$ such that if $(x', 0) \in \mathbb{R}^{2n}$ are the coordinates of $x \in M$, then

$$C_{x_0}(x', t) = (P_{x_0}(x' - x'_0), t) := (u_1, \dots, u_{2n}) = u,$$

where the $(2n-1) \times (2n-1)$ constant matrix P_{x_0} is chosen so that in terms of u -coordinates,

$$(3.2) \quad L_k|_{x_0} = \frac{\partial}{\partial u_{2k-1}} - i\frac{\partial}{\partial u_{2k}}, \quad 1 \leq k \leq n,$$

$$X_0|_{x_0} = -\frac{\partial}{\partial u_{2n-1}}.$$

Note that the second equality of (3.2) actually implies that $X_0|_{(x',0)} = -\partial/\partial u_{2n-1}$ at all points of $M \cap \bar{U}$, and hence that

$$L_n = -i \frac{\partial}{\partial u_{2n-1}} + \frac{\partial}{\partial t}$$

along $M \cap \bar{U}$.

We also note that the matrix P_{x_0} is uniquely determined by the condition (3.2) and uniformly non-singular on \bar{U} , and depends smoothly for all $x_0 \in M \cap \bar{U}$. In terms of u -coordinates, the vector fields L_k , $1 \leq k \leq n-1$, can be written as

$$(3.3) \quad L_k = \left(\frac{\partial}{\partial u_{2k-1}} + \sum_{j=1}^{2n-2} a_j^k(u) \frac{\partial}{\partial u_j} + a_k(u) \frac{\partial}{\partial u_{2n-1}} \right) - i \left(\frac{\partial}{\partial u_{2k}} + \sum_{j=1}^{2n-2} b_j^k(u) \frac{\partial}{\partial u_j} + b_k(u) \frac{\partial}{\partial u_{2n-1}} \right),$$

where a_j^k, b_j^k, a_k, b_k are smooth real valued functions which vanish at 0. In the sequel we let ∂_l , $1 \leq l \leq n$, denote the holomorphic partial derivatives in the l -th variable of local complex valued coordinates. When we change the local coordinate functions, this partial derivative operator will be written in new coordinates. We also let $\tilde{\partial}_\beta$ denotes ∂_β or $\bar{\partial}_\beta$.

Proposition 3.1. *For each $x_0 \in M \cap U$ and positive integer m , there are smooth complex valued coordinates $\zeta = (\zeta_1, \dots, \zeta_n)$, $\zeta_n = t + ix_{2n-1}$, defined near x_0 such that in ζ -coordinates, the vector fields L_1, \dots, L_{n-1} can be written as*

$$(3.4) \quad L_1 = \frac{\partial}{\partial \zeta_1} + \sum_{l=1}^{n-1} a_l(\zeta) \frac{\partial}{\partial \zeta_l} + \sum_{l=1}^{n-1} b_l(\zeta) \frac{\partial}{\partial \bar{\zeta}_l} + \left(e(\zeta) + id(\zeta) \right) \frac{\partial}{\partial x_{2n-1}},$$

$$L_\alpha = \frac{\partial}{\partial \zeta_\alpha} + \sum_{l=1}^{n-1} a_l^\alpha(\zeta) \frac{\partial}{\partial \zeta_l} + \sum_{l=1}^{n-1} b_l^\alpha(\zeta) \frac{\partial}{\partial \bar{\zeta}_l} + \left(e_\alpha(\zeta) + id_\alpha(\zeta) \right) \frac{\partial}{\partial x_{2n-1}},$$

where $2 \leq \alpha \leq n-1$. Also, the coordinate functions satisfy

$$(3.5) \quad \begin{aligned} \partial_1^j \bar{\partial}_1^k b_l(0) &= \partial_1^j \bar{\partial}_1^k a_l^\alpha(0) = \partial_1^j \bar{\partial}_1^k b_l^\alpha(0) = 0, & j+k \leq m, & 2 \leq l \leq n-1, \\ \tilde{\partial}_\beta^i \partial_1^j \bar{\partial}_1^k e(0) &= 0, & i=0,1, & i+j+k \leq m, & 2 \leq \beta \leq n-1, \\ (\partial_1 - \bar{\partial}_1)^s d(0) &= (\partial_1 - \bar{\partial}_1)^s e_\alpha(0) = (\partial_1 - \bar{\partial}_1)^s d_\alpha(0) = 0, & s \leq m. \end{aligned}$$

Proof. Let us take the vector fields L_1, \dots, L_n and smooth coordinates u defined near x_0 so that the vector fields L_k , $1 \leq k \leq n-1$, have the representation as in (3.3). Therefore (3.4) and (3.5) hold for $m = p = 0$.

Assume by induction that there are smooth complex valued coordinates $\zeta = (\zeta_1, \dots, \zeta_n)$ defined near $x_0 \in M$ so that in terms of ζ -coordinates, we can write the vector fields L_ν , $1 \leq \nu \leq n-1$, as in (3.4) where the coefficients satisfy the estimates in (3.5) for m replaced by $p \geq 0$. Set

$$\begin{aligned} x_{2j-1} &= \zeta_j + \bar{\zeta}_j, & x_{2j} &= -i(\zeta_j - \bar{\zeta}_j), & 1 \leq j < n, \\ x_{2n-1} &= -i(\zeta_n - \bar{\zeta}_n), & x_{2n} &= t, \end{aligned}$$

and set

$$D_k = \frac{\partial}{\partial x_k}, \quad 1 \leq k \leq 2n-1.$$

In $x = (x_1, \dots, x_{2n-1}, t)$ coordinates, each vector field L_ν , $1 \leq \nu \leq n-1$, can be written as

$$L_\nu = \left(\frac{\partial}{\partial x_{2\nu-1}} + \sum_{l=1}^{2n-2} c_l^\nu(x) \frac{\partial}{\partial x_l} + e_\nu(x) \frac{\partial}{\partial x_{2n-1}} \right) - i \left(\frac{\partial}{\partial x_{2\nu}} + \sum_{l=1}^{2n-2} d_l^\nu(x) \frac{\partial}{\partial x_l} + d_\nu(x) \frac{\partial}{\partial x_{2n-1}} \right),$$

where $c_l^\nu, d_l^\nu, e_\nu, d_\nu$ are smooth real valued functions. Set

$$\begin{aligned} \tilde{x}_i &= x_i, \quad i = 1, 2, \\ \tilde{x}_l &= x_l - \sum_{\nu=2}^{n-1} \sum_{j+k=p+1} \frac{1}{j!k!} [D_1^j D_2^k c_l^\nu(0) x_1^j x_2^k x_{2\nu-1} + D_1^j D_2^k d_l^\nu(0) x_1^j x_2^k x_{2\nu}] \end{aligned}$$

for $l = 3, \dots, 2n-2$, and set

$$\begin{aligned} \tilde{x}_{2n-1} &= x_{2n-1} - \sum_{j+k=p+1} \frac{1}{(j+1)!k!} [D_1^j D_2^k e_1(0) x_1^{j+1} x_2^k] \\ &\quad - \sum_{\beta=3}^{2n-2} \sum_{j+k=p} \frac{1}{(j+1)!k!} [D_\beta D_1^j D_2^k e_1(0) x_1^{j+1} x_2^k x_\beta] \\ &\quad - \sum_{\nu=2}^{n-1} \left[\frac{1}{(p+1)!} D_2^{p+1} e_\nu(0) x_2^{p+1} x_{2\nu-1} + \frac{1}{(p+1)!} D_2^{p+1} d_\nu(0) x_2^{p+1} x_{2\nu} \right] \\ &\quad - \frac{1}{(p+2)!} D_2^{p+1} d_1(0) x_2^{p+2}. \end{aligned}$$

If we set

$$w_j = \frac{1}{2}(\tilde{x}_{2j-1} + i\tilde{x}_{2j}), \quad 1 \leq j \leq n-1, \quad w_n = \frac{1}{2}(t + i\tilde{x}_{2n-1}),$$

then in w -coordinates, the vector fields L_α , $1 \leq \alpha \leq n-1$, can be written as

$$\begin{aligned} L_1 &= \frac{\partial}{\partial w_1} + \sum_{l=1}^{n-1} \tilde{a}_l(w) \frac{\partial}{\partial w_l} + \sum_{l=1}^{n-1} \tilde{b}_l(w) \frac{\partial}{\partial \bar{w}_l} + (\tilde{e}(w) + i\tilde{d}(w)) \frac{\partial}{\partial \tilde{x}_{2n-1}}, \\ L_\alpha &= \frac{\partial}{\partial w_\alpha} + \sum_{l=1}^{n-1} \tilde{a}_l^\alpha(w) \frac{\partial}{\partial w_l} + \sum_{l=1}^{n-1} \tilde{b}_l^\alpha(w) \frac{\partial}{\partial \bar{w}_l} + (\tilde{e}_\alpha(w) + i\tilde{d}_\alpha(w)) \frac{\partial}{\partial \tilde{x}_{2n-1}}, \end{aligned}$$

where for $2 \leq l, \alpha, \beta \leq n-1$, we have

$$\begin{aligned} \partial_1^j \bar{\partial}_1^k \tilde{b}_l(0) &= 0, \quad \tilde{\partial}_\beta \partial_1^j \bar{\partial}_1^k \tilde{e}(0) = 0, \quad j+k \leq p, \\ \partial_1^j \bar{\partial}_1^k \tilde{a}_l^\alpha(0) &= \partial_1^j \bar{\partial}_1^k \tilde{b}_l^\alpha(0) = \partial_1^j \bar{\partial}_1^k \tilde{e}(0) = 0, \quad j+k \leq p+1, \\ (\partial_1 - \bar{\partial}_1)^s \tilde{d}(0) &= (\partial_1 - \bar{\partial}_1)^s \tilde{e}_\alpha(0) = (\partial_1 - \bar{\partial}_1)^s \tilde{d}_\alpha(0) = 0, \quad s \leq p+1. \end{aligned}$$

We perform the following change of coordinates:

$$\begin{aligned}\zeta_1 &= w_1, \\ \zeta_l &= w_l - \sum_{j+k=p+1} \frac{1}{(j+1)!k!} [\bar{\partial}_1^j \partial_1^k \bar{b}_l(0) \bar{w}_1^{j+1} w_1^k], \quad 2 \leq l \leq n-1.\end{aligned}$$

Then in terms of ζ -coordinates, we may write the vector fields L_1, \dots, L_{n-1} as in (3.4) and the coefficients of L_α , $1 \leq \alpha \leq n-1$, satisfy (3.5) where m is replaced by $p+1$. If we proceed up to m -steps, we will get a proof of the proposition. \square

For each $x_0 \in M \cap U$ and $\delta > 0$, we want to define a quantity $\tau(x_0, \delta)$ in such a way that the successive derivatives of the coefficients of L_k , $1 \leq k \leq n-1$, up to a certain order (less than or equal to m), change by no more than δ on a nonisotropic ball about x_0 . We use the special coordinates $\zeta = (\zeta_1, \zeta_2, \dots, \zeta_n)$ defined near x_0 as in Proposition 3.1.

Assume $x_0 \in M \cap U$. In terms of ζ -coordinates, we write L_k , $1 \leq k \leq n-1$, as in (3.4) such that the coefficient functions satisfy the estimates in (3.5). Note that we may assume that m is an even integer. Let $b(\zeta) = e(\zeta) + id(\zeta)$ be the coefficient function of $\partial/\partial x_{2n-1}$ in L_1 , and set

$$b_{m-1}(\zeta) := 2i \operatorname{Re} \left(\sum_{1 \leq j+k \leq m-1} \frac{1}{j!k!} \partial_1^j \bar{\partial}_1^k d(0) \zeta_1^j \bar{\zeta}_1^k \right) := \sum_{1 \leq j+k \leq m-1} b_{jk}(x_0) \zeta_1^j \bar{\zeta}_1^k.$$

Then, by virtue of the estimates in (3.5), we may write:

$$(3.6) \quad \begin{aligned} b(\zeta) &= \sum_{1 \leq j+k \leq m-1} b_{jk}(x_0) \zeta_1^j \bar{\zeta}_1^k + i \operatorname{Re} \left(\sum_{\beta=2}^{n-1} \sum_{j+k \leq m/2-1} b_{j,k}^\beta(x_0) \zeta_1^j \bar{\zeta}_1^k \bar{\zeta}_\beta \right) \\ &\quad + \mathcal{O}(|\zeta_1|^m + |\zeta'| |\zeta_1|^{\frac{m}{2}} + |\zeta'|^2 + |\zeta_n|), \end{aligned}$$

where $\zeta' = (0, \zeta_2, \dots, \zeta_{n-1}, \zeta_n)$. For $2 \leq \nu \leq n-1$, let us write

$$(3.7) \quad a_\nu(\zeta) = \sum_{1 \leq j+k \leq m/2-1} a_{j,k}^\nu \zeta_1^j \bar{\zeta}_1^k + \mathcal{O}(|\zeta_1|^{\frac{m}{2}} + |\zeta'|),$$

where a_ν 's are coordinate functions of $\partial/\partial \zeta_\nu$ in L_1 . Let $x = (x_1, \dots, x_{2n-1}, t)$ be the real coordinates of ζ and set $D_k = \partial/\partial x_k$, $1 \leq k \leq 2n-1$.

Note that $-ib_{m-1}(\zeta)$ is a smooth real valued function which is an $(m-1)$ -th order polynomial in ζ_1 and $\bar{\zeta}_1$. We let $a(\zeta)$ be a real valued function defined by

$$(3.8) \quad a(\zeta) := \frac{\partial}{\partial x_1} (-ib_{m-1}) = \operatorname{Im} \left[\frac{\partial}{\partial \zeta_1} \bar{b}_{m-1} \right] := \sum_{0 \leq j+k \leq m-2} a_{j,k} \zeta_1^j \bar{\zeta}_1^k.$$

Using the coefficient functions $b_{j,k}^\beta$, $a_{j,k}^\nu$ and $a_{j,k}$ defined in (3.6)-(3.8), we set

$$\begin{aligned} A_l(x_0) &= \max \{ |a_{j,k}| ; j+k = l \}, \quad l = 0, 1, 2, \dots, m-2, \\ E_{l'}(x_0) &= \max \{ |a_{j,k}^\nu|, |b_{j,k}^\beta| ; j+k = l', \beta, \nu = 2, \dots, n-1 \}, \quad 0 \leq l' \leq m/2-1, \end{aligned}$$

and for each $\delta > 0$ we define

$$(3.9) \quad \tau(x_0, \delta) = \min_{\substack{0 \leq l \leq m-2 \\ 0 \leq l' \leq m/2-1}} \{(\delta/A_l(x_0))^{1/(l+2)}, (\delta^{1/2}/E_{l'}(x_0))^{1/(l'+1)}\}.$$

Set $\tau(x_0, \delta) = \tau$ for a convenience. Then it follows from (3.9) that

$$(3.10) \quad |\partial_1^j \bar{\partial}_1^k a(0)| \leq \delta \tau^{-(j+k+2)}, \quad j+k \leq m-2,$$

and for $2 \leq \beta, \nu \leq n-1$, we have

$$(3.11) \quad |\partial_1^j \bar{\partial}_1^k a_\nu(0)| \leq \delta^{1/2} \tau^{-(j+k+1)},$$

$$(3.12) \quad |\bar{\partial}_\beta \partial_1^j \bar{\partial}_1^k b(0)| \leq \delta^{1/2} \tau^{-(j+k+1)}, \quad j+k \leq \text{frac}m2 - 1.$$

If we combine the estimates in (3.10) with the fact that $D_2^s b_{m-1}(0) = 0$, $s \leq m$ (because $(\partial_1 - \bar{\partial}_1)^s d(0) = 0$, $s \leq m$), we obtain, from the expression of $a(\zeta)$ in (3.8), that

$$(3.13) \quad |\tilde{\partial}_\beta^i \partial_1^j \bar{\partial}_1^k b(0)| \leq \delta^{1-i/2} \tau^{-(j+k+1)}, \quad i = 0, 1, \quad i \cdot m/2 + j+k \leq m-1.$$

From the estimates in (3.10)–(3.13) together with the fact that $\eta(L_1) \equiv 0$ it follows, by induction, that

$$(3.14) \quad |\partial_1^j \bar{\partial}_1^k \eta(\frac{\partial}{\partial \zeta_1})(0)| \lesssim \delta^{1/2} \tau^{-(j+k)}, \quad j+k \leq m/2 - 1,$$

because $\eta(\partial/\partial \zeta_\alpha)(0) = 0$, $1 \leq \alpha \leq n-1$, and $\delta^{1/2} \lesssim \tau$.

Set $\omega^n = 1/2(dt + i\eta)$. Since \mathcal{S} is integrable it follows that $\bar{\omega}^n([L_1, L_\alpha]) = 0$ along M and hence we have

$$(3.15) \quad \partial_1^j \bar{\partial}_1^k \eta([L_1, L_\alpha])(0) = 0, \quad j+k \leq m.$$

Combining the estimates in (3.5), (3.10)–(3.15) with the fact that $\eta([L_1, \bar{L}_\alpha])(0) = 0$, $2 \leq \alpha \leq n-1$, one obtains that

$$|\partial_1^j \bar{\partial}_1^k d_\alpha(0)|, |\partial_1^j \bar{\partial}_1^k e_\alpha(0)| \lesssim \delta^{1/2} \tau^{-1}, \quad j+k \leq 1,$$

and hence that

$$|\partial_1^j \bar{\partial}_1^k \eta(\frac{\partial}{\partial \zeta_\alpha})(0)| \lesssim \delta^{1/2} \tau^{-1}, \quad j+k \leq 1.$$

If we use again the estimates in (3.5), (3.10)–(3.15) together with the fact that $\eta(L_\alpha) \equiv 0$, $1 \leq \alpha \leq n-1$, we obtain, by induction, that

$$(3.16) \quad |\partial_1^j \bar{\partial}_1^k \eta(\frac{\partial}{\partial \zeta_\alpha})(0)| \lesssim \delta^{1/2} \tau^{-(j+k)}, \quad j \geq 1, \quad j+k \leq m/2 - 1,$$

$$(3.17) \quad |\partial_1^j \bar{\partial}_1^k d_\alpha(0)|, |\partial_1^j \bar{\partial}_1^k e_\alpha(0)| \lesssim \delta^{1/2} \tau^{-(j+k)}, \quad j \geq 1, \quad j+k \leq m/2 - 1$$

for $2 \leq \alpha \leq n-1$. To obtain the estimates for the derivatives of the form $\bar{\partial}_1^l$ for $l \geq 2$, we note, from the estimates in (3.5), that

$$(3.18) \quad (\bar{\partial}_1 - \partial)^s d_\alpha(0) = (\bar{\partial}_1 - \partial)^s e_\alpha(0) = 0, \quad s \leq m, \quad 2 \leq \alpha \leq n-1.$$

Since we can write

$$(\bar{\partial}_1 - \partial_1)^s = \bar{\partial}_1^s + \sum_{j=1}^{s-1} c_{j,s} \bar{\partial}_1^{s-j} \partial_1^j,$$

where $c_{j,s}$'s are integers, we conclude, from the estimates in (3.16)–(3.18), that

$$(3.19) \quad |\partial_1^j \bar{\partial}_1^k \eta(\frac{\partial}{\partial \zeta_\alpha})(0)| \lesssim \delta^{1/2} \tau^{-(j+k)}, \quad j+k \leq m/2-1,$$

$$(3.20) \quad |\partial_1^j \bar{\partial}_1^k d_\alpha(0)|, |\partial_1^j \bar{\partial}_1^k e_\alpha(0)| \lesssim \delta^{1/2} \tau^{-(j+k)}, \quad j+k \leq m/2-1$$

for $2 \leq \alpha \leq n-1$. By virtue of the estimates in (3.11)–(3.13) together with (3.19) and (3.20), we obtain by induction that

$$(3.21) \quad |\tilde{\partial}_\beta^i \partial_1^j \bar{\partial}_1^k \eta(\frac{\partial}{\partial \zeta_1})(0)| \lesssim \delta^{1-i/2} \tau^{-(j+k+1)}, \quad i=0,1, \quad i \cdot m/2 + j+k \leq m-1.$$

These estimates are essential ingredients to obtain a uniform subelliptic estimates for $\bar{\partial}$ in dilated coordinates.

Set $\tau(x_0, \delta) = \tau$ for convenience and let

$$R_\delta(x_0) = \{\zeta \in \mathbb{C}^n ; |\zeta_1| \leq \tau, |\zeta_\beta| \leq \delta^{1/2} \tau^{-1/2}, \beta=2, \dots, n-1, |\zeta_n| \leq \delta\},$$

and define a dilation map $D_\delta : \mathbb{C}^n \rightarrow \mathbb{C}^n$ by

$$D_\delta(\zeta) = (\tau^{-1} \zeta_1, \delta^{-1/2} \zeta_2, \dots, \delta^{-1/2} \zeta_{n-1}, \delta^{-1} \zeta_n) := (w_1, \dots, w_n) = w.$$

Define

$$Q_\delta(x_0) = D_\delta(R_\delta(x_0)) = \{w \in \mathbb{C}^n ; |w_1| \leq 1, |w_\beta| \leq \tau^{-1/2}, \beta=2, \dots, n-1, |w_n| \leq 1\}$$

and set $L_1^\delta = \tau(D_\delta)_* L_1$. Then, from the expression of L_1 in (3.4), we have

$$\begin{aligned} L_1^\delta &= \frac{\partial}{\partial w_1} + a_1(D_\delta(w)) \frac{\partial}{\partial w_1} + b_1(D_\delta(w)) \frac{\partial}{\partial \bar{w}_1} + \sum_{l=2}^{n-1} \tau \delta^{-1/2} a_l(D_\delta^{-1}(w)) \frac{\partial}{\partial w_l} \\ &\quad + \sum_{l=2}^{n-1} \tau \delta^{-1/2} b_l(D_\delta^{-1}(w)) \frac{\partial}{\partial \bar{w}_l} + \tau \delta^{-1} b(D_\delta^{-1}(w)) \frac{\partial}{\partial y_{2n-1}}, \end{aligned}$$

where $w_n = y_{2n} + iy_{2n-1}$. Set $\eta_\delta = \delta^{-1} \eta$. Therefore, $\eta_\delta(\partial/\partial y_{2n-1}) = 1$ on $M \cap U$.

Set $B(w) = b(D_\delta^{-1}(w))$. Recall the expression of $B(w)$ defined in (3.6). So we can write:

$$\begin{aligned} \tau \delta^{-1} B(w) &= \sum_{1 \leq j+k \leq m-1} b_{jk}(x_0) \delta^{-1} \tau^{j+k+1} w_1 \bar{w}_1^k \\ &\quad + i \operatorname{Re} \left(\sum_{\beta=2}^{n-1} \sum_{j+k \leq m/2-1} b_{j,k}^\beta(x_0) \delta^{-1/2} \tau^{j+k+1} w_1 \bar{w}_1 \bar{w}_\beta \right) + \mathcal{O}(\tau). \end{aligned}$$

Then as in Section 2 of [7], the non-negative condition for $i/2\eta_\delta([L_1^\delta, \bar{L}_1^\delta])$ on $Q_\delta(x_0) \cap M$ forces that

$$(3.22) \quad |\delta^{-1/2}\tau^{j+k+1}b_{j,k}^\beta(x_0)| \lesssim \tau^\gamma \ll 1, \quad j+k \leq \frac{m}{2} - 1,$$

where $\gamma = (10 \times (m/2)!)^{-1}$. Therefore the $b_{j,k}^\beta$'s, in the Taylor expansion of $b(\zeta)$, are not the major terms in the definition of $\tau(x_0, \delta)$ in (3.9), and the estimates in (3.22) show that

$$(3.23) \quad |\bar{\partial}_\beta \partial_1^j \bar{\partial}_1^k b(0)| \lesssim \delta^{1/2} \tau^{-(j+k+1)+\gamma}, \quad j+k \leq m/2 - 1, \quad 2 \leq \beta \leq n-1.$$

By using the estimates in (3.5), (3.14), (3.15), (3.22) and (3.23) together with the fact that $\eta(L_\nu) \equiv 0$, $1 \leq \nu \leq n-1$, we can also show, by the method leading to (3.19) and (3.20), that

$$(3.24) \quad |\partial_1^j \bar{\partial}_1^k \eta(\frac{\partial}{\partial \zeta_\alpha})(0)| \lesssim \delta^{1/2} \tau^{-j-k+\gamma},$$

$$(3.25) \quad |\partial_1^j \bar{\partial}_1^k d_\alpha(0)|, |\partial_1^j \bar{\partial}_1^k e_\alpha(0)| \lesssim \delta^{1/2} \tau^{-j-k+\gamma}, \quad j+k \leq m/2 - 1.$$

Using the estimates in (3.5) and (3.23)–(3.25), we can also show that

$$(3.26) \quad |\partial_1^j \bar{\partial}_1^k \eta([L_1, \bar{L}_\alpha])(0)| \lesssim \delta^{1/2} \tau^{-(j+k+1)+\gamma}, \quad j+k \leq m/2 - 1, \quad 2 \leq \alpha \leq n-1.$$

Now we want to show that the coefficients $a_\nu(\zeta)$ of L_1 satisfy the estimates similar to those in (3.26). Recalling the expression of L_1 in ζ coordinates as in (3.4), we set

$$\tilde{L}_1 = L_1 - \sum_{\nu=2}^{n-1} a_\nu(\zeta) L_\nu, \quad \tilde{L}_\alpha = L_\alpha, \quad 2 \leq \alpha \leq n-1.$$

Then \tilde{L}_1 can be written as:

$$\tilde{L}_1 = \frac{\partial}{\partial \zeta_1} + \sum_{l=1}^{n-1} \tilde{a}_l(\zeta) \frac{\partial}{\partial \zeta_l} + \sum_{l=1}^{n-1} \tilde{b}_l(\zeta) \frac{\partial}{\partial \bar{\zeta}_l} + \tilde{b}(\zeta) \frac{\partial}{\partial x_{2n-1}},$$

where, from the estimates in (3.5), (3.11) and (3.23)–(3.26), we have

$$(3.27) \quad \partial_1^j \bar{\partial}_1^k \tilde{a}_l(0) = \partial_1^j \bar{\partial}_1^k \tilde{b}_l(0) = 0, \quad j+k \leq m-1, \quad 2 \leq l \leq n-1,$$

$$(3.28) \quad |\partial_1^j \bar{\partial}_1^k \tilde{b}(0)| \lesssim \delta \tau^{-(j+k+1)}, \quad j+k \leq m-1,$$

$$(3.29) \quad |\tilde{\partial}_\beta \partial_1^j \bar{\partial}_1^k \tilde{b}(0)| \lesssim \delta^{1/2} \tau^{-(j+k+1)+\gamma}, \quad j+k \leq m/2 - 1$$

for $2 \leq \beta \leq n-1$.

If we combine the estimates in (3.27)–(3.29) and apply the methods leading to the estimates in (3.26), we obtain that

$$(3.30) \quad |\partial_1^j \bar{\partial}_1^k \eta([\tilde{L}_1, \bar{\tilde{L}}_\alpha])(0)| \lesssim \delta^{1/2} \tau^{-(j+k+1)+\gamma}, \quad j+k \leq m/2 - 1, \quad 2 \leq \alpha \leq n-1.$$

Note that

$$\frac{i}{2}\eta([\tilde{L}_1, \bar{L}_\alpha]) = \frac{i}{2}\eta([L_1, \bar{L}_\alpha]) - \frac{i}{2}\sum_{\nu=2}^{n-1} a_\nu(\zeta)\eta([L_\nu, \bar{L}_\alpha]),$$

where

$$\frac{i}{2}\eta([L_\alpha, \bar{L}_\alpha])(0) = c_\alpha(x_0) \geq d_0 > 0, \quad 2 \leq \alpha \leq n-1,$$

for an independent constant $d_0 > 0$. If we use the fact that $i/2\eta([L_\alpha, \bar{L}_\beta])(0) = 0$ for $\alpha \neq \beta$ together with the estimates in (3.11) and (3.26)-(3.30), we obtain that

$$(3.31) \quad |\partial_1^j \bar{\partial}_1^k a_\nu(0)| \lesssim \delta^{1/2} \tau^{-(j+k+1)+\gamma}, \quad j+k \leq m/2-1, \quad 2 \leq \nu \leq n-1.$$

Now set

$$T = L_1 + \bar{L}_1 = \frac{\partial}{\partial x_1} + \sum_{l=1}^{2n-2} a_l(x) \frac{\partial}{\partial x_l} + a_{2n-1}(x) \frac{\partial}{\partial x_{2n-1}}.$$

Then from the estimates of a_{2n-1} in (3.5), and of a_l , $3 \leq l \leq 2n-2$, in (3.31), it follows that

$$(3.32) \quad |\partial_1^j \bar{\partial}_1^k a_l(0)| \lesssim \delta^{1/2} \tau^{-(j+k+1)+\gamma}, \quad j+k \leq m/2-1, \quad 3 \leq l \leq 2n-1.$$

From (3.31), we conclude that the functions $a_\nu(\zeta)$ are also not the major terms in the definition of $\tau(x_0, \delta)$ in (3.9). Therefore we conclude from (3.23) and (3.32) that

$$(3.33) \quad \tau(x_0, \delta) = \min \{(\delta/A_l(x_0))^{1/l+2} ; 0 \leq l \leq m-2\},$$

and hence it follows that $\delta^{1/2} \lesssim \tau \lesssim \delta^{1/m}$, and if $\delta' < \delta''$, then

$$(3.34) \quad (\delta'/\delta'')^{1/2} \tau(x_0, \delta'') \lesssim \tau(x_0, \delta') \lesssim (\delta'/\delta'')^{1/m} \tau(x_0, \delta'').$$

In order to study how $\tau(x_0, \delta)$ depends on x_0 , it is convenient to introduce an analogous quantity $\mu(x, \delta)$ that is defined more intrinsically. Let us cover \bar{M} by a finite number of neighborhoods U_ν , $\nu = 1, \dots, N$, in Ω so that in each U_ν , Proposition 3.1 holds. Let $\{\chi_\nu\}$ be a partition of unity subordinated to the coordinate neighborhoods $\{U_\nu\}$ of Ω , and let m be a given positive integer.

For any $j, k \geq 0$, $j \geq 1$, we define

$$\mathcal{L}_{j,k}^\nu \eta(x) = \frac{i}{2} L_1^{j-1} \bar{L}_1^k \eta([L_1, \bar{L}_1])(x), \quad x \in U_\nu,$$

and set

$$C_l^\nu(x) = \sum_{j+k=l} |\mathcal{L}_{j,k}^\nu \eta(x)|^2, \quad 1 \leq l \leq m-1,$$

and

$$C_l(x) = \sum_{\nu=1}^N \chi_\nu C_l^\nu(x).$$

Set $M = (m + 1)!$, and for each $\delta > 0$, define

$$(3.35) \quad \mu(x, \delta) = \left(\sum_{l=1}^m C_l^{M/l+1}(x) \delta^{-2M/l+1} \right)^{-1/2M}.$$

Note that $\sum_{l=1}^m C_l(x) > 0$ if the type at x is less than or equal to m . Therefore $\mu(x, \delta)$ is defined intrinsically as a smooth function of $\delta > 0$ and x , for x satisfying $\sum_{l=1}^m C_l(x) > 0$. Let us fix $x_0 \in M \cap U$ and take the smooth complex valued coordinates $\zeta = (\zeta_1, \dots, \zeta_n)$ defined on $M \cap U$ as in Proposition 3.1, where

$$\zeta_j = \frac{1}{2}(x_{2j-1} + ix_{2j}), \quad 1 \leq j \leq n-1, \quad \zeta_n = \frac{1}{2}(t + ix_{2n-1}).$$

For each $\delta > 0$, set $\tau_1 = \tau_2 = \tau(x_0, \delta)$ and $\tau_k = \delta^{1/2}$, $3 \leq k \leq 2n-2$, and define

$$(3.36) \quad P_\delta(x_0) = \{x \in \mathbb{R}^{2n} ; |x_i| \leq \tau_i, 1 \leq i \leq 2n-2, |x_{2n-1}| \leq \delta, |t| \leq \delta\}.$$

Without loss of generality we may assume that there is $\nu_0 \in \{1, 2, \dots, N\}$ such that $x_0 \in U_{\nu_0}$ and $\chi_{\nu_0} \geq 1/N$ on $P_\delta(x_0)$. Recall that L_1 can be written, in ζ coordinates, as

$$L_1 = \frac{\partial}{\partial \zeta_1} + \sum_{l=1}^{n-1} a_l(\zeta) \frac{\partial}{\partial \zeta_k} + \sum_{l=1}^{n-1} b_l(\zeta) \frac{\partial}{\partial \bar{\zeta}_l} + b(\zeta) \frac{\partial}{\partial x_{2n-1}},$$

where, from the estimates in (3.5), (3.23) and (3.31), we have

$$(3.37) \quad \begin{aligned} |\partial_1^j \bar{\partial}_1^k a_l(0)|, |\tilde{\partial}_\beta^j \partial_1^j \bar{\partial}_1^k b(0)| &\lesssim \delta^{1/2} \tau^{-(j+k+1)+\gamma}, \quad j+k \leq m/2-1, \\ |\partial_1^j \bar{\partial}_1^k b(0)| &\lesssim \delta \tau^{-(j+k+1)}, \quad j+k \leq m-1, \\ \partial_1^j \bar{\partial}_1^k b_l(0) &= 0, \quad j+k \leq m, \quad 2 \leq l, \beta \leq n-1. \end{aligned}$$

Therefore may write

$$(3.38) \quad \mathcal{L}_{j,k} \eta(\zeta) = -\partial_1^{j-1} \bar{\partial}_1^k \left[\operatorname{Im} \left(\frac{\partial}{\partial \zeta_1} \bar{b}(\zeta) \right) \eta \left(\frac{\partial}{\partial x_{2n-1}} \right) \right] + E_{j+k},$$

where $E_{j,k}$ satisfies, from the estimates in (3.37), that

$$(3.42) \quad |\tilde{\partial}_\beta^i \partial_1^{l_1} \bar{\partial}_1^{l_2} E_{j+k}(0)| \lesssim \delta^{1-i/2} \tau^{-(j+k+l_1+l_2+1)+\gamma}$$

for $i = 0, 1$, and $i \cdot m/2 + j + k + l_1 + l_2 \leq m - 1$. From (3.37)-(3.39) it follows that

$$(3.40) \quad |\tilde{\partial}_\beta^i \partial_1^{l_1} \bar{\partial}_1^{l_2} \mathcal{L}_{j,k} \eta(0)| \lesssim \delta^{1-i/2} \tau^{-(j+k+l_1+l_2+1)},$$

for $i = 0, 1$, and $i \cdot m/2 + j + k + l_1 + l_2 \leq m - 1$. By (3.39), (3.40) combined with a simple Taylor's theorem argument, we then have

$$(3.41) \quad |\mathcal{L}_{j,k} \eta(\zeta)| \lesssim \delta \tau^{-(j+k+1)}, \quad \zeta \in P_\delta(x_0).$$

By virtue of the definition of $\mu(x, \delta)$ in (3.35), (3.41) implies that

$$(3.42) \quad \mu(x, \delta) \gtrsim \tau(x_0, \delta), \quad x \in P_\delta(x_0).$$

Conversely, let us show that $\mu(x, \delta) \lesssim \tau(x_0, \delta)$ for $x \in P_\delta(x_0)$. Recall that $\tau(x_0, \delta)$ is actually defined as in (3.33). Set

$$(3.43) \quad T(x_0, \delta) = \min\{l ; (\delta/A_l(x_0))^{1/l+2} = \tau(x_0, \delta)\}.$$

Therefore there must exist integers j_0, k_0 with $(j_0 - 1) + k_0 = T(x_0, \delta)$ such that

$$(3.44) \quad \left| \frac{1}{(j_0 - 1)!k_0!} \partial_1^{j_0-1} \bar{\partial}_1^{k_0} a(0) \right| = \left| \frac{1}{(j_0 - 1)!k_0!} \partial_1^{j_0-1} \bar{\partial}_1^{k_0} \left[\operatorname{Im} \frac{\partial}{\partial \zeta_1} \bar{b} \right](0) \right| = \delta \tau^{-j_0-k_0-1}.$$

Assuming that $\chi_{\nu_0} \geq 1/N$, we obtain from the estimates in (3.38), (3.39) and (3.44) that

$$|\mathcal{L}_{j_0, k_0}^{\nu_0} \eta(x_0)| \geq \frac{1}{2N} (j_0 - 1)! k_0! \delta \tau(x_0, \delta)^{-j_0-k_0-1},$$

provided that δ is sufficiently small. Again, by using the estimates in (3.38)–(3.44) and the Taylor series method, we obtain that

$$(3.45) \quad |\mathcal{L}_{j_0, k_0}^{\nu_0} \eta(x)| \approx \delta \tau(x_0, \delta)^{-j_0-k_0-1}, \quad x \in P_\delta(x_0).$$

If we combine (3.45) and the definition of $\mu(x, \delta)$, we obtain that

$$(3.46) \quad \mu(x, \delta) \lesssim \tau(x_0, \delta), \quad x \in P_\delta(x_0).$$

Combining (3.42) and (3.46), we have proved the following proposition.

Proposition 3.2. *If $x \in P_\delta(x_0)$, then*

$$(3.47) \quad \tau(x_0, \delta) \approx \mu(x, \delta).$$

Corollary 3.3. *Suppose $x \in P_\delta(x_0)$. Then*

$$(3.48) \quad \tau(x_0, \delta) \approx \tau(x, \delta)$$

Proof. If we set $x = x_0$ in (3.47), we see that $\mu(x, \delta) \approx \tau(x_0, \delta)$. Since this holds for $x_0 = x$, it follows that $\mu(x, \delta) \approx \tau(x, \delta)$. Hence (3.48) follows. \square

Remark 3.4. $\mu(x, \delta)$ is defined intrinsically. That is, it does not depend on the choice of a specific coordinates. Propositions 3.2 and Corollary 3.3 show that the quantity $\tau(x, \delta)$ is also defined invariantly, up to a universal constant, with respect to the coordinate functions.

Now we want to construct a family of plurisubharmonic functions with large Hessian. The existence of these functions will be a crucial ingredient in the subelliptic estimates for $\bar{\partial}$ -type equation. Note that we are free to choose $x_0 \in M$ and $\delta > 0$. Now assume that $x^\nu \in M$. Let us take the special coordinates $\zeta = (\zeta_1, \dots, \zeta_n)$ defined near x^ν and write the vector fields L_1, \dots, L_n as in (3.4) satisfying (3.5). Also, let $T(x_0, \delta)$ be defined in (3.43). Let $x = (x', t)$ be the real coordinates for ζ , where $(x^\nu, 0) = x^\nu \in M$. Set $\tau_1 = \tau(x^\nu, \delta)$, $\tau_\beta = \delta^{1/2}$, $2 \leq \beta \leq n-1$, and $\tau_n = \delta$, and put $\alpha = (\alpha_1, \dots, \alpha_n)$.

Proposition 3.5. *Suppose $x^\nu \in M \cap U$. Then there exist a small constant $a > 0$ and a smooth function $h_{x^\nu, \delta}$ on $M \cap U$ which satisfy the following:*

- (i) $|h_{x^\nu, \delta}(x)| \leq 1$ and $h_{x^\nu, \delta} \in C_0^\infty(P_\delta(x^\nu))$.
- (ii) If $|t| \leq a\delta$ and if $h_{x^\nu, \delta}$ is not plurisubharmonic at $x = (x', t)$, then

$$T(x', a\delta) < T(x^\nu, \delta).$$

- (iii) If $x \in P_{a\delta}(x^\nu)$, $|t| \leq a\delta$, and if the inequality

$$\partial\bar{\partial}h_{x^\nu, \delta}(x)(L, \bar{L}) \gtrsim \prod_{k=1}^n \tau_k^{-2} |b_k|^2$$

fails to hold at $x = (x', t)$ for $L = \sum_{j=1}^n b_j L_j$, then $T(x', a\delta) < T(x^\nu, \delta)$.

- (iv) For all $x \in P_\delta(x^\nu)$ and all $L = \sum_{j=1}^n b_j L_j$ at x ,

$$|\partial\bar{\partial}h_{x^\nu, \delta}(L, \bar{L})| \lesssim \prod_{k=1}^n \tau_k^{-2} |b_k|^2.$$

- (v) $|D^\alpha h_{x^\nu, \delta}(x)| \leq C_\alpha \prod_{k=1}^n \tau_k^{-\alpha_k}$, where $D^\alpha = \partial_1^{\beta_1} \bar{\partial}_1^{\gamma_1} \dots \partial_n^{\beta_n} \bar{\partial}_n^{\gamma_n}$, $\alpha_i = \beta_i + \gamma_i$.

Proof. We use the definitions of $\mu(x, \delta)$, $\tau(x, \delta)$, $T(x, \delta)$ and the fact that $\tau(x^\nu, \delta) \approx \mu(x, \delta)$ for $x \in P_\delta(x^\nu)$. Then the proof of (i)–(v) follow the same lines as in the proof of Proposition 3.1 in [7]. \square

For each $\varepsilon > 0$, we set $\Omega_\varepsilon = M \times (-1, \varepsilon)$ and $S(\varepsilon) = M \times (-\varepsilon, \varepsilon)$. By adding up the functions $h_{x^\nu, \delta}$ constructed in Proposition 3.5, we can construct bounded plurisubharmonic weight functions so that the Hessian of these functions satisfy certain essentially maximal bounds in a thin strip $S(\varepsilon)$ of M_0 . The heart of these construction is the so-called “doubling property” of $P_\delta(x_0)$, which comes from the relation in (3.48). For a detailed proof of the following theorem, one can refer to Section 3 of [4]. For each small $\delta > 0$, we set $\tau_1(x) = \tau(x, \delta)$, $\tau_2(x) = \dots = \tau_{n-1}(x) = \delta^{1/2}$ and $\tau_n(x) = \delta$ as before.

Theorem 3.6. *For all small $\delta > 0$, there is a plurisubharmonic function $\psi_\delta \in C^\infty(\Omega_\delta)$ which satisfies*

- (i) $|\psi_\delta(x)| \leq 1$, $x \in U \cap \Omega_\delta$.
- (ii) For all $L = \sum_{j=1}^n b_j L_j$ at $x \in U \cap S(\delta)$,

$$(3.49) \quad \partial\bar{\partial}\psi_\delta(x)(L, \bar{L}) \approx \sum_{j=1}^n |b_j(x)|^2 \tau_j^{-2}(x), \quad \text{and}$$

- (iii) $|D^\alpha \psi_\delta(x)| \lesssim C_\alpha \prod_{k=1}^n \tau_k^{-\alpha_k}(x)$, where $D^\alpha = \partial_1^{\beta_1} \bar{\partial}_1^{\gamma_1} \dots \partial_n^{\beta_n} \bar{\partial}_n^{\gamma_n}$, $\alpha_i = \beta_i + \gamma_i$.

In D_2 -equation, we will assign a Dirichlet condition on one side of bS_g^+ , and the Neumann condition on the other side of bS_g^+ . This fact leads us to another difficulty which was not occurred in 1/2-subelliptic estimates of Catlin in [5]. To overcome this difficulty, we need the following Lemma 3.7 and Proposition 3.8 which will be used in the proof of subelliptic estimates for D_2 -equation in Section 5.

Lemma 3.7. *Let $x^\nu \in M \cap U$ and set $c_{1k}^n = \omega^n([L_1, \bar{L}_k])$, $1 \leq k \leq n$. Then for each small $\delta > 0$, we have*

$$(3.50) \quad |c_{11}^n(x)| \lesssim \delta \tau(x^\nu, \delta)^{-2}, \quad x \in P_\delta(x^\nu),$$

$$(3.51) \quad |c_{1k}^n(x)| \lesssim \delta^{1/2} \tau(x^\nu, \delta)^{-1+\gamma}, \quad x \in P_\delta(x^\nu), \quad 2 \leq k \leq n,$$

where $\gamma = (10 \times (m/2)!)^{-1}$.

Proof. Along $M \cap U$, we have $dt = 1/2(dt + i\eta) + 1/2(dt - i\eta)$ and $L_n = \partial/\partial t - i\partial/\partial x_{2n-1}$, which imply that $\partial t = 1/2(dt + i\eta)$. Hence (3.50) follows from (3.41). Since \mathcal{L}^0 is integrable to infinite order along M_0 , it follows that

$$(3.52) \quad \partial \bar{\partial} t(L_1, \bar{L}_n) = \partial \bar{\partial} t(L_1 + \bar{L}_1, \bar{L}_n) = \frac{i}{2} \eta \left(\left[L_1 + \bar{L}_1, \frac{\partial}{\partial t} + i \frac{\partial}{\partial x_{2n-1}} \right] \right)$$

along $M \cap U$. Note that we can write $L_1 + \bar{L}_1 = T = T_1 + tT_2 + x_{2n-1}T_3$, where the coefficient functions of T_1 does not depend on t or x_{2n-1} . From (3.5) together with the estimates of the coefficient functions in (3.37) it follows that

$$(3.53) \quad \left| \partial_1^j \bar{\partial}_1^k \frac{\partial}{\partial t} \eta(T_1)(0) \right|, \quad \left| \partial_1^j \bar{\partial}_1^k \frac{\partial}{\partial x_{2n-1}} \eta(T_1)(0) \right| \lesssim \delta^{1/2} \tau^{-(j+k+1)+\gamma}, \quad j+k \leq m/2 - 1.$$

Since $\eta(L_1 + \bar{L}_1) = \eta(T) \equiv 0$, we have

$$(3.54) \quad \partial_1^j \bar{\partial}_1^k \frac{\partial}{\partial t} \eta(T)(0) = \partial_1^j \bar{\partial}_1^k \frac{\partial}{\partial x_{2n-1}} \eta(T)(0) = 0, \quad j+k \leq m-1.$$

Combining (3.53) and (3.54), we obtain that

$$|\partial_1^j \bar{\partial}_1^k \eta(T_2)(0)|, \quad |\partial_1^j \bar{\partial}_1^k \eta(T_3)(0)| \lesssim \delta^{1/2} \tau^{-(j+k+1)+\gamma}, \quad j+k \leq m/2 - 1.$$

Therefore it follows that

$$|\partial_1^j \bar{\partial}_1^k \eta([T, \frac{\partial}{\partial t} + i \frac{\partial}{\partial x_{2n-1}}])(0)| \lesssim \delta^{1/2} \tau^{-(j+k+1)+\gamma}, \quad j+k \leq m/2 - 1,$$

and this proves (3.51) for $k = n$.

When $2 \leq k \leq n-1$, we use the estimates in (3.26). □

For each $\delta > 0$, let ψ_δ be the function constructed in Theorem 3.6. We need the following proposition which will be used to prove Lemma 5.6 that is necessary for the estimates of (5.39) and (5.40) in the subelliptic estimates of D_2 operator in Section 5.

Proposition 3.8. *For all small $\delta > 0$ and for each $\alpha = (\alpha_1, \dots, \alpha_n)$ we have*

$$(3.54) \quad |c_{11}^n(x) D^\alpha \psi_\delta(x)| \leq C_\alpha \delta \tau(x, \delta)^{-2} \prod_{k=1}^n \tau_k(x, \delta)^{-\alpha_k},$$

$$(3.55) \quad |c_{1,l}^n(x) D^\alpha \psi_\delta(x)| \leq C_\alpha \delta^{\frac{1}{2}} \tau(x, \delta)^{-1+\gamma} \prod_{k=1}^n \tau_k(x, \delta)^{-\alpha_k}$$

for $x \in S(\delta) \cap U$. Here $\tau_1(x, \delta) = \tau(x, \delta)$, $\tau_k(x, \delta) = \delta^{1/2}$, $2 \leq k \leq n-1$ and $\tau_n(x, \delta) = \delta$.

Proof. Note that the functions ψ_δ in Theorem 3.6 were constructed by adding up functions $h_{x^\nu, \delta}$, $x^\nu \in M \cap U$, constructed in Proposition 3.5, with $\text{supp } h_{x^\nu, \delta} \subset P_\delta(x^\nu)$. By virtue of (3.48), there is a small $c > 0$, independent of $\delta > 0$, so that we can arrange points $x^\nu = x^\nu(\delta) \in M \cap U$, $\nu \in I$, satisfying

$$S(\delta) \cap U \subset \cup_{\nu \in I} P_\delta(x^\nu), \quad \text{and} \quad P_{c\delta}(x^\nu) \cap P_{c\delta}(x^\mu) = \emptyset \quad \text{if} \quad \nu \neq \mu.$$

Then, as in the proof of Lemma 3.3 in [4], there is a fixed integer N (independent of δ) such that any $(N+1)$ intersection of $P_\delta(x^\nu)$'s are empty.

Now assume that $x \in S(\delta) \cap U$. Then there are $\nu_1, \dots, \nu_l \in I$, $l \leq N$, such that $x \in P_\delta(x^{\nu_j})$, $1 \leq j \leq l$. By virtue of (3.48) again, it follows that

$$\tau(x, \delta) \approx \tau(x^{\nu_1}, \delta) \approx \dots \approx \tau(x^{\nu_l}, \delta),$$

independent of δ . If we express the vector fields L_1, \dots, L_n in terms of the special coordinates in each neighborhood of x^{ν_i} , $i = 1, \dots, l$, then (3.54) and (3.55) follow from (3.50) and (3.51), respectively. \square

4. SPECIAL FRAMES FOR ALMOST COMPLEX STRUCTURES

Assume $\overline{M} \subset \widetilde{M}$ and let $\varphi \in C^\infty(\overline{M})$ be a smooth real-valued function such that $\varphi(x) > 0$ for $x \in M$, and $\varphi(x) = 0$, $d\varphi(x) \neq 0$ for $x \in bM$. We can extend φ to Ω by requiring that it be independent of t . Let us denote by T_p the type at a point $p \in \overline{M}$ and define

$$T(\overline{M}) = \max\{T_p ; p \in \overline{M}\}.$$

Since the type condition is an open condition, it follows that $T(\overline{M})$ is well-defined and is finite. In the sequel, we assume that $T(\overline{M}) = m < \infty$. We define $r \in C^\infty(\Omega)$ by $r(x, t) = t(\phi(x))^{-2m}$ and for any ε, σ , $0 < \varepsilon \leq \sigma \leq 1$, define

$$(4.1) \quad S_{\varepsilon, \sigma} = \{(x, t) \in \Omega ; \varphi(x) > 0 \text{ and } 0 \leq r(x, t) \leq \varepsilon \sigma^{2m}\}.$$

Remark 4.1. The quantities ε and σ will be fixed later. If we set $g(x) = \varepsilon \cdot \sigma^{2m} \cdot \varphi(x)^{2m}$, then g is the required positive function in the definition of S_g^+ in Section 1 and $S_{\varepsilon, \sigma}$ equals to S_g^+ .

We define a subbundle of \mathcal{L}^0 on $S_{\varepsilon, \sigma}$ by letting $\mathcal{R}_{(x, t)} = \{L \in \mathcal{L}_{(x, t)}^0 ; Lr = 0\}$. Clearly, the map H defined by $H(L) = L - (Lr)(L_n r)^{-1} L_n$ defines an isomorphism of \mathcal{S} onto \mathcal{R} (at all points of $S_{\varepsilon, \sigma}$). Set $\mu_1(x) = \mu(x, \varepsilon \phi(x)^{2m})$, $\mu_2(x) = \dots$, $\mu_{n-1}(x) = \varepsilon^{1/2} \phi(x)^m$, $\mu_n(x) = \varepsilon \phi(x)^{2m}$. We define a weighted metric $\langle \cdot, \cdot \rangle$ on \mathcal{L}^0 by the relations:

$$\begin{aligned} \langle H(L_j), H(L_k) \rangle &= \mu_j(x)^{-1} \mu_k(x)^{-1} \langle L_j, L_k \rangle_0, & 1 \leq j, k \leq n-1 \\ \langle L_n, L_n \rangle &= \varepsilon^{-2} \varphi(x)^{-4m}, \\ \langle L_n, H(L_l) \rangle &= 0, & 1 \leq l \leq n-1, \end{aligned}$$

where $L_l \in \mathcal{S}$, $1 \leq l \leq n-1$. Since $\mu(x, \delta)$ is a smooth function of x and δ , it follows that $\langle \cdot, \cdot \rangle$ is a smooth Hermitian metric on \mathcal{L}^0 . Now, using Proposition 3.1, we shall cover $S_{\varepsilon, \sigma}$ by special (dilated) coordinate neighborhoods such that on each of them, there is a frame \mathcal{L} that satisfying required good estimates.

Proposition 4.2. *There exist constants ε_0 and σ_0 such that if $0 < \varepsilon < \varepsilon_0$ and $0 < \sigma < \sigma_0$, then on $S_{\varepsilon, \sigma}$ there exists for all $x_0 \in M$ with $\varphi(x_0) > 0$ a neighborhood $W(x_0) \subset S_{\varepsilon, \sigma}$ with the following properties:*

- (i) *On $W(x_0)$ there are smooth coordinates y_1, \dots, y_{2n} so that $W(x_0) = \{y ; |y'| < \sigma, 0 \leq y_{2n} \leq \sigma^{2m}\}$, where $y' = (y_1, y_2, \dots, y_{2n-1})$ is independent of t and the function y_{2n} is defined by $y_{2n} = \varepsilon^{-1} \varphi(x)^{-2m} t$. Thus, $M_0 \cap W(x_0)$ and $M_\sigma \cap W(x_0)$ correspond to the points in $W(x_0)$ where $y_{2n} = 0$ and σ^{2m} , respectively. Moreover, the point $(x_0, 0) \in \Omega$, which we identify with x_0 , corresponds to the origin.*
- (ii) *The above coordinate charts are uniformly smoothly related in the sense that if $W(\tilde{p}_0)$ and $W(x_0)$ intersect, and if \tilde{y} and y_0 are the associated coordinates, then*

$$|D^\alpha(\tilde{y} \circ (y_0)^{-1})| \leq C_{|\alpha|}$$

holds on that portion of \mathbb{R}^{2n} where $\tilde{y} \circ (y_0)^{-1}$ is defined. The constant $C_{|\alpha|}$ is independent of ε , σ , and x_0 .

- (iii) *On $W(x_0)$, there exists a smooth frame $\{L_1, \dots, L_n\}$ for \mathcal{L} such that if $\{\omega^1, \dots, \omega^n\}$ is the dual frame, and if L_k and ω^k are written as $\sum_{j=1}^{2n} b_{kj} \partial / \partial y_j$ and $\sum_{j=1}^{2n} d_{kj} dy_j$, then*

$$\sup_{y \in W(x_0)} \left\{ |D_y^\alpha b_{kj}(y)| + |D_y^\alpha d_{kj}(y)| \right\} \leq C_{|\alpha|},$$

where $C_{|\alpha|}$ is independent of x_0 , j, k, ε and σ .

- (iv) *With the frames as in (iii), set $c_{1l}^n = \omega^n([L_1, \bar{L}_l])$, $l = 2, \dots, n$. Then there is a constant $C > 0$ independent of x_0 , ε and σ such that*

$$(4.2) \quad \sup_{y \in W(x_0)} |c_{1l}^n(y)| \leq C \sigma^{2m}.$$

- (v) *There are constants $c > 0$ and $C > 0$ independent of x_0 , ε and σ such that if $B_b(x)$ denotes the ball of radius b about $x \in S_{\varepsilon, \sigma}$ with respect to the metric $\langle \cdot, \cdot \rangle$, then*

$$(4.3) \quad B_{c\sigma}(x_0) \subset W(x_0) \subset B_{C\sigma}(x_0),$$

and if $\text{Vol } B_b(x_0)$ denotes the volume of $B_b(x_0)$ with respect to $\langle \cdot, \cdot \rangle$, then

$$(4.4) \quad cb^{2n-1} \sigma^{2m} \leq \text{Vol } B_b(x_0) \leq Cb^{2n-1} \sigma^{2m}.$$

Proof. We first cover \bar{M} by a finite number of neighborhoods V_ν , $\nu = 1, \dots, N$, in Ω such that in each V_ν there exist coordinates (u_1, \dots, u_{2n}) with the property that $u_{2n} = t$ and that $u_k(u', t) = u_k(u', 0)$, $k < 2n$, for $(u', t) \in V_\nu$, and that $\partial / \partial u_{2n-1} = -X_0$ at all points of $M \cap V_\nu$. Also, we can arrange the neighborhoods V_ν so that Proposition 3.1 holds on each V_ν .

For any point $x_0 \in M \cap V_\nu$, we take coordinate functions $\zeta^\nu = (\zeta_1^\nu, \dots, \zeta_n^\nu)$ constructed as in Proposition 3.1. Let us set $\zeta^\nu = \zeta$ and denote by L_k^ν , the vector fields L_k , $1 \leq k \leq n$, written in ζ^ν -coordinates, and let $x = (x_1, \dots, x_{2n-1}, t)$ be the real coordinates of ζ . In ζ -coordinates, we may write:

$$L_k^\nu = \frac{\partial}{\partial \zeta_k} + \sum_{l=1}^{n-1} a_l^k(\zeta) \frac{\partial}{\partial \zeta_l} + \sum_{l=1}^{n-1} b_l^k(\zeta) \frac{\partial}{\partial \bar{\zeta}_l} + \left(e_k(\zeta) + id_k(\zeta) \right) \frac{\partial}{\partial x_{2n-1}}, \quad 1 \leq k \leq n,$$

where the coefficients a_l^k, b_l^k, e_k, d_k satisfy the estimates in (3.5) and (3.37).

Set $\delta = \varepsilon\varphi(x_0)^{2m}$, and let us write $b(\zeta) = e_1(\zeta) + id_1(\zeta)$, where $e_1(\zeta)$ and $d_1(\zeta)$ are smooth real valued functions. Recall the Taylor expansion of $b(\zeta)$ at x_0 (i.e., at $\zeta = 0$) in (3.6). We then take the quantity $\tau(x_0, \delta)$ and the corresponding quantity $\mu(x, \delta)$ defined in (3.9) and (3.35), respectively. By Proposition 3.5, it follows that $\mu(x_0, \delta) \approx \tau(x_0, \delta)$, and hence the estimates in (3.37) imply that

$$(4.5) \quad |\partial_1^j \bar{\partial}_1^k b(0)| \lesssim \delta \mu(x_0, \delta)^{-j-k-1}, \quad j+k \leq m-1,$$

$$(4.6) \quad |\tilde{\partial}_\beta \partial_1^j \bar{\partial}_1^k b(0)| \lesssim \delta^{1/2} \mu(x_0, \delta)^{-(j+k+1)+\gamma}, \quad j+k \leq m/2-1,$$

where $\tilde{\partial}_\beta = \partial_\beta$ or $\bar{\partial}_\beta$, $2 \leq \beta \leq n-1$, and $\partial_k = \partial/\partial\zeta_k$, $1 \leq k \leq n$, and $\gamma = (10 \times (m/2)!)^{-1}$.

Set $\mu = \mu(x_0, \delta)$ for a convenience. We define new coordinates $y = D_{\varepsilon, x_0}(x) = (y_1, \dots, y_{2n})$ by means of dilation map $D_{\varepsilon, x_0} : \mathbb{R}^{2n} \rightarrow \mathbb{R}^{2n}$ given by

$$(4.7) \quad y = (\mu^{-1}x_1, \mu^{-1}x_2, \delta^{-1/2}x_3, \dots, \delta^{-1/2}x_{2n-2}, \delta^{-1}x_{2n-1}, e^{-1}\varphi(x)^{-2m}x_{2n}),$$

where $\varphi(x)$ is the function φ expressed in the x -coordinates of x_0 . In terms of the y -coordinates we define an open set $W_b(x_0)$ by

$$(4.8) \quad W_b(x_0) = \{x \in V_\nu \cap S_{\varepsilon, \sigma} ; |y_k(x)| < b, 1 \leq k \leq 2n-1, 0 \leq y_{2n}(x) \leq \sigma^{2m}\}.$$

Note that in $W_b(x_0)$, $y_{2n} = 0$ and $y_{2n} = \sigma^{2m}$ coincide with $r = 0$ and $r = \varepsilon\sigma^{2m}$, respectively, on the boundaries of $S_{\varepsilon, \sigma}$. We define a frame $\{L_1, \dots, L_n\}$ in $W_b(x_0)$ by setting

$$(4.9) \quad \begin{aligned} L_1 &= \mu(x, \delta)(L_1^\nu - r_1 L_n^\nu) = \mu(x, \delta)H(L_1^\nu), \\ L_k &= \varepsilon^{1/2}\varphi(x)^m(L_k^\nu - r_k L_n^\nu) = \varepsilon^{1/2}\varphi(x)^m H(L_k^\nu), \quad 2 \leq k \leq n-1, \\ L_n &= \varepsilon\varphi(x)^{2m}L_n^\nu, \end{aligned}$$

where $r_k = (L_k^\nu r)(L_n^\nu r)^{-1}$, and L_k^ν is the vector field L_k written in x -coordinates of V_ν . Set

$$L_k^\nu = \sum_{l=1}^{2n-1} e_{kl}(x) \frac{\partial}{\partial x_l}, \quad 1 \leq k \leq n-1, \quad \text{and} \quad L_n^\nu = \frac{\partial}{\partial t} + \sum_{l=1}^{2n-1} e_{nl}(x) \frac{\partial}{\partial x_l}.$$

In terms of dilated coordinates y in $W_b(x_0)$, we set $E_{kl} = e_{kl} \cdot D_{\varepsilon, x_0}^{-1}$, $R_k = r_k \cdot D_{\varepsilon, x_0}^{-1}$, $\Phi = \varphi \cdot D_{\varepsilon, x_0}^{-1}$ and $\Phi_l = (\partial\varphi/\partial x_l) \cdot D_{\varepsilon, x_0}^{-1}$.

By a direct calculation, one obtains that

$$(4.10) \quad R_k = \frac{2\varepsilon m \sum_{l=1}^{2n-1} E_{\nu, l} \Phi^{2m-1} \Phi_l y_{2n}}{1 + 2\varepsilon m \sum_{l=1}^{2n-1} E_{n, l} \Phi^{2n-1} \Phi_l y_{2n}}, \quad 1 \leq k \leq n-1.$$

We set

$$|f|_{k, W_b(x_0)} = \sup \{|D_y^\alpha f(y)| ; y \in W_b(x_0), |\alpha| \leq k\},$$

and extend this norm to vector fields and 1-forms by using the coefficients of $\partial/\partial y_j$ or dy_j . By virtue of Proposition 3.2, it follows that for all $x \in W_b(x_0)$, we have

$$(4.11) \quad \mu(x, \delta) \approx \tau(x_0, \delta) \lesssim \delta^{1/m} = \varepsilon^{1/m} \varphi(x_0)^2.$$

Combining (4.10) with (4.11), we conclude that for each $s \geq 0$ there are C_s independent of x_0 and δ such that

$$(4.12) \quad |\delta^{-1}\mu(x, \delta)R_k|_{s, W_b(x_0)} \leq C_s b \varepsilon^{1/m}, \quad 1 \leq k \leq n-1.$$

Therefore L_1 can be written, in y -coordinates, as

$$(4.13) \quad L_1 = \frac{\mu(x, \delta)}{\mu(x_0, \delta)} \sum_{l=1}^2 E_{1,l} \frac{\partial}{\partial y_l} + \mu(x, \delta) \varepsilon^{-1/2} \varphi(x_0)^{-m} \sum_{l=3}^{2n-2} E_{1,l} \frac{\partial}{\partial y_l} \\ + \mu(x, \delta) \varepsilon^{-1} \varphi(x_0)^{-2m} E_{1,2n-1} \frac{\partial}{\partial y_{2n-1}} + \mathcal{O}(\tilde{E}),$$

where

$$(4.14) \quad |E_{1,1} - 1|_{s, W_b(x_0)} \leq C \mu(x_0, \delta), \quad s \leq m, \\ |E_{1,2} + i|_{s, W_b(x_0)} \leq C \mu(x_0, \delta), \quad s \leq m,$$

and \tilde{E} satisfies, from the estimates in (4.12), that

$$(4.15) \quad |\tilde{E}|_{s, W_b(x_0)} \leq C \varepsilon^{1/m}, \quad s \leq m,$$

for an independent constant $C > 0$. By virtue of the estimates in (3.31), we also have

$$(4.16) \quad |\mu(x, \delta) \varepsilon^{-1/2} \varphi(x_0)^{-m} E_{1,l}|_{k, W_b(x_0)} \leq C \tau(x_0, \delta)^\gamma, \quad k \leq m+1, \quad 3 \leq l \leq 2n-2.$$

Observe that the diameter in the x -coordinates of $W_b(x_0)$ is $\mathcal{O}(b\mu(x_0, \delta)) \ll \varphi(x_0)$ by (4.11). Hence it is clear that $\mu(x, \delta)\mu(x_0, \delta)^{-1}$ and $\Phi\varphi(x_0)^{-1}$ are very close to 1 (independent of x_0 and δ) in $W_b(x_0)$ if b is small. We also observe that $\mu(x, \delta)$ is defined independently of coordinate functions and that $\tau(x_0, \delta) \approx \tau(x, \delta) \approx \mu(x, \delta)$ if $x \in W_b(x_0)$. Therefore it follows that $\mu(x_0, \delta) \approx \mu(p_0, \delta)$ and $\varphi(x_0) \approx \varphi(p_0)$ if $W_b(x_0) \cap W_b(p_0) \neq \emptyset$. These facts prove (ii).

Now set

$$w_k = \frac{1}{2}(y_{2k-1} - iy_{2k}), \quad 1 \leq k \leq n-1, \quad \text{and} \quad w_n = \frac{1}{2}(y_{2n} - iy_{2n-1}),$$

and define $D_l = \partial/\partial y_l$, $1 \leq l \leq 2n$, and $B_{m-1}(y) = b_{m-1} \circ D_{\varepsilon, x_0}^{-1}(y)$, where b_{m-1} is the $(m-1)$ -th order polynomial of $b(\zeta)$ in ζ_1 and $\bar{\zeta}_1$ as in (3.6). We recall that the real part $e(\zeta)$ of the coefficient function of $\partial/\partial x_{2n-1}$ in L_1 satisfies the estimates in (3.5). Hence, by combining (4.12)–(4.16), we obtain that

$$(4.17) \quad L_1 = \frac{\partial}{\partial w_1} + \mu(x_0, \delta) \delta^{-1} B_{m-1}(y) \frac{\partial}{\partial y_{2n-1}} + \mathcal{O}(E),$$

$$(4.18) \quad X = L_1 + \bar{L}_1 = \frac{\partial}{\partial y_1} + \mathcal{O}(E),$$

where

$$(4.19) \quad |E|_{s, W_b(x_0)} \leq C_s \varepsilon^{\gamma/m}, \quad |\mu(x_0, \delta) \delta^{-1} B_{m-1}|_{s, W_b(x_0)} \leq C_s$$

for an independent constant $C > 0$. Combining (4.17)–(4.19), we conclude that if $b \leq \sqrt{\sigma}$,

$$\lim_{\sigma \rightarrow 0} \left| L_1 - \left(\frac{\partial}{\partial w_1} + \mu(x_0, \delta) \delta^{-1} B_{m-1}(y) \frac{\partial}{\partial y_{2n-1}} \right) \right|_{s, W_b(x_0)} = 0, \quad s \leq m+1.$$

Setting $W(x_0) = W_\sigma(x_0)$ for sufficiently small σ , we obtain (i) and (iii).

To prove (iv), we recall that $L_n = \partial/\partial y_{2n} - i\partial/\partial y_{2n-1} + y_{2n}T$, where $Ty_{2n} \equiv 0$, and that \mathcal{L} is integrable to infinite order along M_0 . Hence

$$(4.20) \quad c_{1n}^n = \omega^n([L_1, \bar{L}_n]) = \omega^n \left(\left[L_1 + \bar{L}_1, \frac{\partial}{\partial y_{2n}} - i \frac{\partial}{\partial y_{2n-1}} \right] \right) + \mathcal{O}(y_{2n}).$$

Combining (4.18) – (4.20), we conclude that

$$(4.21) \quad \sup_{y \in W(x_0)} |c_{1n}^n(y)| \leq C_1(\varepsilon^{\gamma/m} + \sigma^{2m}),$$

where $\gamma = (10 \times (m/2)!)^{-1}$.

For the estimates of c_{1l}^n , $l = 2, \dots, n-1$, we need the estimates in (3.26). Then it follows that $|c_{1l}^n(y)| \leq C_1 \varepsilon^{\gamma/m}$, along M_0 . Since L_1 and L_l are tangential vector fields, for $2 \leq l \leq n-1$, it follows that

$$(4.22) \quad |c_{1l}^n(y)| \leq C_1(\varepsilon^{\gamma/m} + |y_{2n}|) \leq C_1(\varepsilon^{\gamma/m} + \sigma^{2m}), \quad y \in W(x_0),$$

Now we assume that $\varepsilon \leq \sigma^{2m^2 \cdot \gamma^{-1}}$. Then (4.2) follows from (4.21) and (4.22).

By Proposition 3.5, it follows that $\tau(x_0, \delta) \approx \mu(x, \delta)$ for $x \in W(x_0)$. Since $\{L_1, \dots, L_n\}$ is orthonormal with respect to $\langle \cdot, \cdot \rangle$, we conclude that (4.3) and (4.4) hold if σ is sufficiently small. \square

Using the special coordinates y_1, \dots, y_{2n} and the special frames L_1, \dots, L_n defined in (4.9), we want to define L^2 -operators with mixed boundary conditions. In the process of subelliptic estimates for D_2 -operator, we will see that certain boundary integral terms on M_0 occur. To handle these boundary integral terms, we need the following lemma.

Lemma 4.3. *There are a frame $\{X_1, \dots, X_n\}$ for \mathcal{L} and its dual frame $\{\eta^1, \dots, \eta^n\}$ so that if we set $c_{kn}^n = \eta^n([X_k, \bar{X}_n])$, $1 \leq k \leq n-1$, then*

$$(4.23) \quad \begin{aligned} c_{kn}^n &= 0 \quad \text{on } W(x_0), \quad k = 2, \dots, n-1, \\ |c_{1n}^n|_{s, W(x_0)} &\leq C_s \sigma^{2m}. \end{aligned}$$

Proof. With the frames L_1, \dots, L_n and $\omega^1, \dots, \omega^n$ defined on $W(x_0)$ in (4.9), we let $b_k(y)$, $2 \leq k \leq n-1$, be the smooth function satisfying

$$(4.24) \quad \omega^n([L_k, \bar{L}_n])(y) + \sum_{l=2}^{n-1} \bar{b}_l(y) \omega^n([L_k, \bar{L}_l])(y) \equiv 0, \quad y \in W(x_0).$$

Since the Levi-form of L_2, \dots, L_{n-1} is always positive definite, (4.24) is solvable on $W(x_0)$. Set

$$X_k = L_k, \quad k = 1, \dots, n-1, \quad \text{and} \quad X_n = L_n + \sum_{l=2}^{n-1} b_l(y) L_l.$$

Then its dual frames η^k , $k = 1, \dots, n$, are given by

$$\eta^k = \omega^k - b_k(y)\omega^n, \quad k = 1, \dots, n-1, \quad \text{and} \quad \eta^n = \omega^n.$$

In terms of the new frames it follows from (4.24) that

$$c_{kn}^n := \eta^n([X_k, \bar{X}_n]) = \omega^n([L_k, \bar{L}_n]) + \sum_{l=2}^{n-1} \bar{b}_l \omega^n([L_k, \bar{L}_n]) \equiv 0, \quad 2 \leq k \leq n-1,$$

on $W(x_0)$ and

$$(4.25) \quad c_{1n}^n := \omega^n([X_1, \bar{X}_n]) = \omega^n([L_1, \bar{L}_n]) + \sum_{l=2}^{n-1} \bar{b}_l(y)\omega^n([L_1, \bar{L}_l]).$$

Therefore (4.23) follows from (4.25) and (4.2). \square

Recall that a deformation of \mathcal{L}^0 is a section A of the bundle $\Gamma^1(S_{\varepsilon, \sigma})$. In terms of the special frames in $W(x_0)$, we write $A = \sum_{j,l=1}^n A_{jl} \bar{\omega}^l \otimes L_j$, and then define

$$|A(y)|_s = \sum_{|\alpha| \leq s} \sum_{j,l=1}^n |D_y^\alpha A_{jl}(y)|,$$

$$|A|_{s, W(x_0)} = \sup\{|A(y)|_s ; y \in W(x_0)\}.$$

We suppose that A satisfies

$$(4.26) \quad |A|_{m+2n+3, W(x_0)} \leq \varepsilon_0$$

for a sufficiently small $\varepsilon_0 > 0$.

We define $\mathcal{A}(S_{\varepsilon, \sigma})$ to be the space of sections $A \in \Gamma^{0,1}(S_{\varepsilon, \sigma} ; 0)$ such that along M_0 , $A(\bar{L}) = 0$ whenever $\bar{L} \in T^{0,1} \cap \mathcal{CT}M_0$. From now on, we assume that $A \in \mathcal{A}(S_{\varepsilon, \sigma})$. Then we can define a deformation \mathcal{L}^A of \mathcal{L}^0 by

$$\bar{\mathcal{L}}^A = \{\bar{L} + A(\bar{L}) ; \bar{L} \in \bar{\mathcal{L}}_z^0, \quad z \in S_{\varepsilon, \sigma}\}.$$

In terms of the frames X_1, \dots, X_n , and its dual frames η^1, \dots, η^n in $W(x_0)$ constructed in Lemma 4.3, we define

$$X_j^A = X_j + \bar{A}(X_j), \quad j = 1, \dots, n,$$

and let η_A^j be the dual frames. Set

$$(4.27) \quad L_j^A = \sigma^{1/4} \left(X_j^A - (X_j^A r)(X_n^A r)^{-1} X_n^A \right), \quad 1 \leq j \leq n-1, \quad L_n^A = X_n^A$$

and

$$(4.28) \quad \omega_A^j = \sigma^{-1/4} \eta_A^j, \quad 1 \leq j \leq n-1, \quad \omega_A^n = \left(\eta_A^n + \sum_{j=1}^{n-1} (X_j^A r)(X_n^A r)^{-1} \eta_A^j \right).$$

Obviously, the frames ω_A^j , for $j = 1, \dots, n$, are dual to L_j^A , and $L_j^A r \equiv 0$ for $1 \leq j \leq n-1$.

Note that we can write $L_n = \partial/\partial y_{2n} - i\partial/\partial y_{2n-1} + y_{2n}T$, where $Ty_{2n} \equiv 0$. Assuming that A satisfies (4.26) for sufficiently small $\varepsilon_0 \leq \sigma^{2m^2 \cdot \gamma^{-1}}$, it follows from Lemma 4.3 that

$$(4.29) \quad \sup_{y \in W(x_0)} |\omega_A^n([L_k^A, \bar{L}_n^A])(y)| \leq C\sigma^{1/4}|A|_1, \quad 2 \leq k \leq n-1,$$

$$(4.30) \quad \sup_{y \in W(x_0)} |\omega_A^n([L_1^A, \bar{L}_n^A])(y)| \leq C\sigma^{2m+1/4},$$

where the constant $C > 0$ is independent of x_0, σ and ε .

In order to measure how L_j^A , $j = 1, 2, \dots, n$, depend on A , we define

$$P_k(y; A) = \sum_{\substack{k_1, \dots, k_N \\ |k_1| + \dots + |k_N| \leq k}} \prod_{\nu=1}^N |A(y)|_{k_\nu}.$$

In the sequel, we assume that A satisfies (4.26) for sufficiently small ε_0 .

Lemma 4.4.. *For $y \in W(x_0)$, the following pointwise estimates hold:*

$$(4.31) \quad |L_k^A - \sigma^{1/4}L_k|_s \leq C_s P_s(y; A), \quad |\omega_A^k - \sigma^{-1/4}\omega^k|_s \leq C_s \sigma^{-1/4} P_s(y; A)$$

for $1 \leq k \leq n-1$ and

$$(4.32) \quad |L_n^A - L_n|_s \leq C_s P_s(y; A), \quad |\omega_A^n - \omega^n|_s \leq C_s P_s(y; A).$$

Proof. From the expression of L_k^A and ω_A^k in (4.27) and (4.28), the error terms are the finite product of derivatives as in (4.31) and (4.32). \square

For the subelliptic estimates on the non-euclidean balls $W(x_0)$, we still have to construct a family of plurisubharmonic functions with maximal Hessian in dilated coordinates y defined in (4.7). By virtue of Theorem 3.6 there is a family of plurisubharmonic functions $\{\psi_\delta(x)\}_{\delta>0}$ defined on $\Omega_\delta \cap U = \{(x', t) ; t \leq \delta\} \cap U$. We may assume that there is an open set $\widetilde{W}(x_0) = W_{C\sigma}(x_0)$, for some $C > 1$, such that $W(x_0) = W_\sigma(x_0) \Subset \widetilde{W}(x_0) \Subset P_{a\delta}(x_0)$, provided σ is sufficiently small. We define

$$S_\rho = \{y \in \Omega_\rho \cap U ; |y_{2n}| < \rho\},$$

and set $\delta_0 = \delta(x_0) = \varepsilon\phi(x_0)^{2m}$ and

$$(4.33) \quad \begin{aligned} \mu_1(x) &= \mu(x, \delta_0), \quad \mu_k(x) = \varepsilon^{1/2}\varphi(x)^m, \quad 2 \leq k \leq n-1, \\ \mu_n(x) &= \varepsilon\varphi(x)^{2m}. \end{aligned}$$

For any $\rho > 0$ we set

$$(4.34) \quad \mu_1^\rho(x) = \mu(x, \rho\delta_0), \quad \mu_k^\rho(x) = (\rho\delta_0)^{1/2}, \quad 2 \leq k \leq n-1, \quad \mu_n^\rho(x) = \rho\delta_0.$$

Theorem 4.5. *For each small $\rho > 0$, there exists a C^∞ plurisubharmonic function λ_ρ defined on $\widetilde{W}(x_0) \ni W(x_0)$ satisfying the following:*

- (i) $|\lambda_\rho| \leq 1$ in $\widetilde{W}(x_0)$.
- (ii) For all $y \in S_\rho \cap \widetilde{W}(x_0)$, and $L^A = \sum_{j=1}^n b_j L_j^A$, we have

$$\partial\bar{\partial}\lambda_\rho(y)(L^A, \bar{L}^A) \approx \sum_{k=1}^n |b_k|^2 \mu_k^2 (\mu_k^\rho)^{-2}.$$

$$(iii) |L^A \lambda_\rho|^2 \lesssim \partial\bar{\partial}\lambda_\rho(y)(L^A, \bar{L}^A).$$

$$(iv) |D^\alpha \lambda_\rho| \leq C_\alpha \prod_{k=1}^n \mu_k^{\alpha_i}(x) (\mu_k^\rho)^{-\alpha_i}, \text{ where } D^\alpha = \partial_1^{\beta_1} \bar{\partial}_1^{\gamma_1} \dots \partial_n^{\beta_n} \bar{\partial}_n^{\gamma_n}, \alpha_i = \beta_i + \gamma_i.$$

Proof. Let $\{\psi_{\rho\delta_0}\}_{\rho\delta_0>0}$ be the family of plurisubharmonic functions constructed in Theorem 3.6. Set $\lambda_\rho(y) = \psi_{\rho\delta_0} \circ D_{x_0}^{-1}(y)$, where D_{x_0} is the dilation function defined in (4.8). It is clear that λ_ρ is plurisubharmonic and satisfies (i). Note that the orthonormal frame $\{L_1, \dots, L_n\}$ defined in (4.9) can be written as $L_j = \mu_j(x)(L_j^\nu - r_j L_n^\nu)$, $1 \leq j \leq n-1$, and $L_n = \mu_n(x)L_n^\nu$.

If we let $L = \sum_{j=1}^n b_j L_j$, then it follows by functoriality that

$$(4.35) \quad \begin{aligned} \partial\bar{\partial}\lambda_\rho(y)(L, \bar{L}) &= \partial\bar{\partial}\psi_{\rho\delta_0}(x)(dD_{x_0}^{-1}L, dD_{x_0}^{-1}\bar{L}) \\ &= \partial\bar{\partial}\psi_{\rho\delta_0}(x) \left(\sum_{j=1}^n b_j \mu_j L_j^\nu, \sum_{k=1}^n \bar{b}_k \mu_k \bar{L}_k^\nu \right) \\ &\quad - 2\operatorname{Re} \left[\sum_{j=1}^{n-1} \sum_{k=1}^n r_j b_j \bar{b}_k \mu_j \mu_k \partial\bar{\partial}\psi_{\rho\delta_0}(x)(L_n^\nu, \bar{L}_k^\nu) \right] \\ &\quad + \sum_{j,k=1}^{n-1} r_j \bar{r}_k b_j \bar{b}_k \mu_j \mu_k \partial\bar{\partial}\psi_{\rho\delta_0}(x)(L_n^\nu, \bar{L}_n^\nu). \end{aligned}$$

From the expression of $r_k(x)$ in (4.10) and (4.12) it follows that

$$(4.36) \quad |r_k(x)| \lesssim \varepsilon^{1/m} (\rho\delta_0) \cdot \mu(x, \delta_0)^{-1} \quad \text{for } |t| \leq \rho\delta_0.$$

If we combine (4.36) with the properties (ii) and (iii) of Theorem 3.6, we conclude that

$$(4.37) \quad \partial\bar{\partial}\lambda_\rho(y)(L, \bar{L}) \approx \sum_{k=1}^n |b_k|^2 \mu_k^2 (\mu_k^\rho)^{-2}.$$

Note that the vector fields L_j^A 's and its dual frames ω_j^A 's were written in terms of L_j 's, and η_A^j 's as in (4.27) and (4.28). By virtue of the expressions in (4.31) and (4.32) we can write

$$(4.38) \quad \partial\bar{\partial}\lambda_\rho(y)(L^A, \bar{L}^A) = \partial\bar{\partial}\lambda_\rho(y)(L, \bar{L}) + \mathcal{O}(A^2 \rho^{-2} + A \rho^{-3/2}) |L|^2.$$

Recall that $A = 0$ along M and satisfies (4.26). Hence $|A|_0 \lesssim \varepsilon_0 \rho^2$ provided $|y_{2n}| \leq \rho$. Therefore it follows from (4.37) and (4.38) that

$$\partial\bar{\partial}\lambda_\rho(y)(L^A, \bar{L}^A) \approx \sum_{k=1}^n |b_k|^2 \mu_k(x)^2 (\mu_k^\rho)^{-2}.$$

This proves (ii). (iii) follows from the estimate

$$|L^A \lambda_\rho|^2 \lesssim \sum_{k=1}^n |b_k|^2 |L_k^A \lambda_\rho|^2 \lesssim \sum_{k=1}^n |b_k|^2 \mu_k(x)^2 (\mu_k^\rho)^{-2} \approx \partial \bar{\partial} \lambda_\rho(y)(L^A, \bar{L}^A).$$

(iv) follows from the property (iii) of Theorem 3.6. \square

Next, we show that there exists a smooth Hermitian metric on $S_{\varepsilon, \sigma}$ such that for all $x_0 \in M$ the frames L_1^A, \dots, L_n^A given by (4.27) are orthonormal. For $L \in \mathcal{L}^0$ and $A \in \mathcal{A}(S_{\varepsilon, \sigma})$ satisfying (4.26), define a bundle isomorphism $P_A : \mathcal{L}^0 \rightarrow \mathcal{L}^A$ by $P_A(L) = L + A(L)$. Define a homomorphism $H_A : \mathcal{L}^A \rightarrow \mathcal{R}^A$, where $\mathcal{R}^A = \{L \in \mathcal{L}^A ; Lr = 0\}$, by

$$H_A(L) = L - \frac{Lr}{X_n^A r} X_n^A = L - \frac{Ly_{2n}}{L_n^A y_{2n}} L_n^A.$$

Then $H_A \circ P_A$ is an isomorphism of \mathcal{R} onto \mathcal{R}^A . We define a metric $\langle \cdot, \cdot \rangle_A$ on \mathcal{L}^A by

$$\begin{aligned} \langle (H_A \circ P_A) \bar{L}_1, (H_A \circ P_A) \bar{L}_2 \rangle_A &= \langle \bar{L}_1, \bar{L}_2 \rangle, \quad \bar{L}_1, \bar{L}_2 \in \mathcal{R}, \\ \langle L_n^A, L_n^A \rangle_A &= 1, \\ \langle (H_A \circ P_A) \bar{L}_1, L_n^A \rangle_A &= 0, \quad \bar{L}_1 \in \mathcal{R}. \end{aligned}$$

Note that L_n^A is actually globally defined, so that the above conditions determine a metric on \mathcal{L}^A . Since L_j , $j = 1, 2, \dots, n-1$, defined in (4.9), are an orthonormal basis of \mathcal{L} , it follows that $L_j^A = (H_A \circ P_A)L_j$, $j = 1, 2, \dots, n-1$, are an orthonormal basis of \mathcal{L}^A with respect to $\langle \cdot, \cdot \rangle_A$.

Let dV denote the volume form associated with the Riemannian metric $\langle \cdot, \cdot \rangle$. In the coordinates (y_1, \dots, y_{2n}) in $W(x_0)$ we can write $dV = V(y)dy$, where $dy = dy_1 \dots dy_{2n}$, and where V satisfies

$$|V|_{s, W(x_0)} \leq C_s \quad \text{and} \quad \inf_{y \in W(x_0)} V(y) > c > 0,$$

where c is independent of σ , ε , and x_0 . We will define the inner product for two functions $g, h \in C^\infty(S_{\varepsilon, \sigma})$ by

$$(g, h) = \int g \bar{h} dV.$$

Let $\Lambda^{0, q}(S_{\varepsilon, \sigma}; A)$ denote the space of $(0, q)$ -forms with respect to \mathcal{L}^A on $S_{\varepsilon, \sigma}$, and set

$$\Gamma^{0, q}(S_{\varepsilon, \sigma}; A) = \Lambda^{0, q}(S_{\varepsilon, \sigma}; A) \otimes \mathcal{L}^A.$$

Now let us define, for a given structure \mathcal{L}^A , where A satisfying (4.26) for small ε_0 , the L^2 -operators corresponding to $D_q : \Gamma^{0, q} \rightarrow \Gamma^{0, q+1}$ and its adjoint. We define $\mathcal{E}_c^{0, q}(S_{\varepsilon, \sigma}; A)$ to be the set of smooth sections U of $\Gamma^{0, q}(S_{\varepsilon, \sigma}; A)$ such that support of U is a compact subset of $S_{\varepsilon, \sigma}$. Let $\mathcal{E}_0^{0, q}(S_{\varepsilon, \sigma}; A)$ denote the set of sections of $\mathcal{E}_c^{0, q}(S_{\varepsilon, \sigma}; A)$ with compact support in the interior of $S_{\varepsilon, \sigma}$. Suppose that $U = \sum_{\ell=1}^n \sum_{|J|=q} U_\ell^J \bar{\omega}_A^J \cdot L_\ell^A$ is an element of $\Gamma^{0, q}(S_{\varepsilon, \sigma}; A)$ with compact support in $W(x_0)$. We define

$$(4.39) \quad \|U\|^2 = \int_{S_{\varepsilon, \sigma}} \sum_{\ell=1}^n \sum_{|J|=q} |U_\ell^J|^2 dV,$$

where dV is the volume form given by the metric of \mathcal{L}^0 . Since $\{L_1^A, \dots, L_n^A\}$ is an orthonormal frame, the quantity in (4.39) is independent of the frame neighborhood $W(x_0)$. Thus, by using a partition of unity, it follows that the norm in (4.39) extends to all of $\Gamma^{0,q}(S_{\varepsilon,\sigma}; A)$. Let $L_q^2(S_{\varepsilon,\sigma}, T_A^{1,0})$ denote the set of sections of $\Gamma^{0,q}(S_{\varepsilon,\sigma}; A)$ such that (4.39) is finite.

Define $\mathcal{B}_-^q(S_{\varepsilon,\sigma}; A)$ to be the set of forms in $\mathcal{E}_c^{0,q}(S_{\varepsilon,\sigma}; A)$ such that U_ℓ^J vanishes on M_0 whenever $n \notin J$. (This is also independent of the frame neighborhood $W(x_0)$.) Similarly, define $\mathcal{B}_+^q(S_{\varepsilon,\sigma}; A)$ to be the set of forms in $\mathcal{E}_c^{0,q}(S_{\varepsilon,\sigma}; A)$ such that U_ℓ^J vanishes on M_σ whenever $n \in J$. We now define the formal adjoint D'_q of D_q on $\mathcal{E}_c^{0,q}(S_{\varepsilon,\sigma}; A)$ by $D'_q U = G \in \mathcal{E}_c^{0,q-1}(S_{\varepsilon,\sigma}; A)$ if for all $V \in \mathcal{E}_0^{0,q-1}(S_{\varepsilon,\sigma}; A)$,

$$(U, D_q V) = (G, V),$$

where $(\ , \)$ corresponds to the norm in (4.39). By combining (2.3) with (2.6) together with integration by parts, it follows that if $U = \sum_\nu U_\nu L_\nu^A \in \Gamma^{0,k}(S_\varepsilon; A)$ is supported in $W(x_0)$, then

$$(4.40) \quad D'_k U = \sum_\nu \left(\bar{\partial}^* U_\nu - \sum_\mu \sum_j \partial \bar{\omega}_A^\mu(L_j^A, \bar{L}_\nu^A)(\bar{L}_j^A \lrcorner U_\mu) \right) L_\nu^A,$$

where

$$(4.41) \quad \begin{aligned} \bar{\partial}^* U_\nu = & - \sum_{|J|=k-1} \sum_{j=1}^n (L_j^A U_\nu^{jJ} + e_j U_\nu^{jJ}) \bar{\omega}_A^J \\ & - \sum_{|K|=k-2} \sum_{\ell=1}^n \sum_{i < j} \omega_A^\ell([L_i^A, L_j^A]) U_\nu^{ijK} \bar{\omega}_A^{\ell K}. \end{aligned}$$

We now extend the definition of the operator D_q and D'_q to the L^2 -spaces. Define $L_k^2(S_{\varepsilon,\sigma}; A)$ to be the set of all sections U of $\Gamma^{0,k}(S_{\varepsilon,\sigma}; A)$ for which $\|U\|^2 < \infty$. We define an operator

$$T : L_{q-1}^2(S_{\varepsilon,\sigma}; T_A^{1,0}) \rightarrow L_q^2(S_{\varepsilon,\sigma}; T_A^{1,0})$$

by the condition that $U \in \text{Dom}(T)$ and $TU = F \in L_q^2(S_{\varepsilon,\sigma}; T_A^{1,0})$ if for all $V \in \mathcal{B}_-^q(S_{\varepsilon,\sigma}; A)$, we have

$$(U, D'_q V) = (F, V).$$

Similarly, we can define $S : L_q^2(S_{\varepsilon,\sigma}; T_A^{1,0}) \rightarrow L_{q+1}^2(S_{\varepsilon,\sigma}; T_A^{1,0})$. Note that these definitions imply that if $U \in \text{Dom}(T)$ (or $\text{Dom}(S)$), then $TU = D_q U$ (or $SU = D_{q+1} U$) as in the sense of distribution theory. Let $T^* : L_q^2(S_{\varepsilon,\sigma}; T_A^{1,0}) \rightarrow L_{q-1}^2(S_{\varepsilon,\sigma}; T_A^{1,0})$ and $S^* : L_{q+1}^2(S_{\varepsilon,\sigma}; T_A^{1,0}) \rightarrow L_q^2(S_{\varepsilon,\sigma}; T_A^{1,0})$ be the Hilbert space adjoints of T and S , respectively. It follows that if $U \in \text{Dom}(T^*)$ and $V \in \text{Dom}(S^*)$, then

$$(4.42) \quad T^* U = D'_q(U) \quad \text{and} \quad S^* V = D'_{q+1}(V),$$

as in the sense of distributions. Therefore it follows that

$$\begin{aligned} \mathcal{E}_c^{0,q-1}(S_{\varepsilon,\sigma}; A) \cap \text{Dom}(T) &= \mathcal{B}_+^{q-1}(S_{\varepsilon,\sigma}; A), \\ \mathcal{E}_c^{0,q}(S_{\varepsilon,\sigma}; A) \cap \text{Dom}(T^*) &= \mathcal{B}_-^q(S_{\varepsilon,\sigma}; A). \end{aligned}$$

Similar relations hold for S . Set

$$\mathcal{B}^q(S_{\varepsilon,\sigma}; A) = \mathcal{B}_+^q(S_{\varepsilon,\sigma}; A) \cap \mathcal{B}_-^q(S_{\varepsilon,\sigma}; A).$$

Then we can approximate $U \in \text{Dom}(S) \cap \text{Dom}(T^*)$ by $U_\mu \in \mathcal{B}^q(S_{\varepsilon,\sigma}; A)$ in the graph norm of S and T^* [5, Lemma 6.4].

Lemma 4.6. *Let $U \in \text{Dom}(S) \cap \text{Dom}(T^*)$. Then there exists $U_\mu \in \mathcal{B}^q(S_{\varepsilon,\sigma}; A)$ such that*

$$\lim_{\mu \rightarrow \infty} (\|U_\mu - U\| + \|SU_\mu - SU\| + \|T^*U_\mu - T^*U\|) = 0.$$

Finally, suppose that we have proved the estimate

$$(4.43) \quad \|U\|^2 \leq C(\|T^*U\|^2 + \|SU\|^2)$$

for all $U \in \mathcal{B}^q(S_{\varepsilon,\sigma}; A)$. Then Lemma 4.6 shows that (4.43) holds for all $U \in \text{Dom } T^* \cap \text{Dom } S$. Then from the usual $\bar{\partial}$ -Neumann theory it follows that for all $G \in L_q^2(S_{\varepsilon,\sigma}; T_A^{1,0})$, there exists an element $NG \in \text{Dom}(T^*) \cap \text{Dom}(S)$ such that

$$\|NG\| \leq C^2 \|G\|,$$

and

$$(G, V) = (T^*(NG), T^*V) + (SNG, SV), \quad V \in \text{Dom}(T^*) \cap \text{Dom}(S).$$

We will call N the Neumann operator associated with D_q .

5. THE SUBELLIPTIC ESTIMATE FOR D_q

In this section we prove a subelliptic estimate for the D_q -Neumann problem with almost complex structure \mathcal{L}^A . We set $q = 2$ in this section.

We first define tangential norms that will be used in the estimates. For any $s \in \mathbb{R}$, set

$$\|f\|_s^2 = \int_0^{\sigma^{2m}} \int_{\mathbb{R}^{2n-1}} |\hat{f}(\xi, y_{2n})|^2 (1 + |\xi|^2)^s d\xi dy_{2n},$$

where $\hat{f}(\xi, y_{2n}) = \int_{\mathbb{R}^{2n-1}} e^{-iy' \cdot \xi} f(y', y_{2n}) dy'$. For any integer $k \geq 0$ and any $s \in \mathbb{R}$, set

$$\|f\|_{s,k}^2 = \sum_{j=0}^k \left\| \left\| \frac{\partial^j f}{\partial y_{2n}^j} \right\| \right\|_{s-j}^2.$$

Then, for any integer $m \geq 0$ and $f \in C^\infty(W(x'))$, set

$$\|f\|_m^2 = \sum_{|\alpha| \leq m} \|D_y^\alpha f\|^2.$$

By using the coefficients of U , we can easily define all of the above norms for any section U of $\Gamma^{0,q}$. We recall that $\mathcal{A}(S_{\varepsilon,\sigma})$ is the space of sections $A \in \Gamma^{0,1}(S_{\varepsilon,\sigma}; 0)$ such that along M_0 , $A(\bar{L}) = 0$ whenever $\bar{L} \in T^{0,1} \cap \text{CTM}_0$. We let $C > 1$ and $0 < c \leq 1$ be independent constants which may vary in various estimations. Then the goal of this section is to prove the following subelliptic estimate:

Theorem 5.1. *Suppose $T(\overline{M}) = m < \infty$ and that A is a section of $\mathcal{A}(S_{\varepsilon,\sigma})$. Then there exist small positive constants σ_1 and ε_1 so that if $\varepsilon < \varepsilon_1$, $\sigma < \sigma_1$, and $|A|_{m+2n+3, W(x_0)} \leq \varepsilon$, then the D_q -Neumann problem on $S_{\varepsilon,\sigma}$ with coefficient σ for the almost complex structure \mathcal{L}^A satisfies the following estimate for all forms $U \in \mathcal{B}^q(S_{\varepsilon,\sigma}; A)$ that are compactly supported in $W(x_0)$*

$$(5.1) \quad \sigma^{-3} \|U\|^2 + L^A(U) + \sigma^{\frac{1}{2}} \|U\|_{\frac{1}{m}, 1}^2 \leq C(\|SU\|^2 + \|T^*U\|^2),$$

where $L^A(U)$ is defined by

$$(5.2) \quad \begin{aligned} L^A(U) = & \sum_{k=1}^{n-1} \sum_{\ell=1}^n \sum_{|J|=q} \|\overline{L}_k^A U_\ell^J\|^2 + \sum_{k=1}^{n-1} \sum_{\ell=1}^n \sum_{\substack{|J|=q \\ n \notin J}} \|L_k^A U_\ell^J\|^2 \\ & + \sum_{\ell=1}^n \sum_{\substack{|J|=q \\ n \in J}} \|L_n^A U_\ell^J\|^2 + \sum_{\ell=1}^n \sum_{\substack{|J|=q \\ n \notin J}} \|\overline{L}_n^A U_\ell^J\|^2. \end{aligned}$$

We first state some necessary lemmas for the proof of Theorem 5.1 [5, Lemma 7.5].

Lemma 5.2. *Let $X_j = \sum_k a_{jk} \partial/\partial x_k$, $j = 1, \dots, \ell$, be smooth compactly supported vector fields in \mathbb{R}^d , and suppose that there exist a set $K \Subset \mathbb{R}^d$ and a constant $c > 0$ such that for all $x \in K$,*

$$c < \inf \left\{ \sum_j |\eta(X_j)| + \sum_{i < j} |\eta([X_i, X_j])| ; \eta \in T_x^*, |\eta| = 1 \right\}.$$

Then there exists a constant C independent of X_1, \dots, X_ℓ such that for all $u \in C_0^\infty(\mathbb{R}^d)$ with $\text{supp } u \subset K$ and any integer $s > (d+5)/2$,

$$c \|u\|_{1/2}^2 \leq C \left[\sum_{j=1}^{\ell} \|X_j u\|^2 + \|u\|^2 \sum_{j,k} \|a_{j,k}\|_s^2 \right].$$

If $n \geq 3$, then the $(n-2)$ positive eigenvalue condition on the Levi-form of M guarantees the existence of at least one positive eigenvalue. Set $X_{2k-1} = \text{Re} L_k^A$ and $X_{2k} = \text{Im} L_k^A$ for $1 \leq k \leq n-1$. Then Lemma 5.2 and the expression of L_k^A in (4.27) show that

Lemma 5.3. *Assume that $n \geq 3$. Then for all $f \in C_0^\infty(W(x_0))$,*

$$(5.3) \quad \sigma^{1/2} \|f\|_{1/2}^2 \leq C \sum_{k=1}^{n-1} \left(\|L_k^A f\|^2 + \|\overline{L}_k^A f\|^2 \right) + C \|f\|^2.$$

For convenience, in what that follows we omit the notation A from the frames L_1^A, \dots, L_n^A , and $\omega_A^1, \dots, \omega_A^n$, and L^A . Note that in $W(x_0)$, we have technically chosen coordinates in such a way that $y_{2n} = 0$ and $y_{2n} = \sigma^{2m}$ coincide with $r = 0$ and $r = \varepsilon \sigma^{2m}$, respectively, on the boundaries of $S_{\varepsilon,\sigma}$. Then the following lemma can be proved by modifying the proof of Lemma 7.7 in [5].

Lemma 5.4. *Suppose that $f \in C_0^\infty(W(x_0))$ and f vanishes either on M_0 or on M_σ . If σ is sufficiently small, say $\sigma < \sigma_1$, then there exists a constant C independent of ε , σ , and x_0 such that*

$$(5.4) \quad \sigma^{-4} \|f\|^2 \leq C(\|f\|_{1/m}^2 + \|\tilde{L}_n f\|^2),$$

where $\tilde{L}_n = L_n$ or \bar{L}_n .

To handle the commutator terms, we need the following lemma.

Lemma 5.5. *Assume that $n \geq 3$. Let $U \in \mathcal{B}_q(S_{\varepsilon, \sigma}; A)$ be compactly supported in $W(x_0)$ and assume that $|K| = q - 1$ with $n \notin K$ and that $1 \leq k \leq n - 1$. Set $c_{kn}^n = w^n([L_k, \bar{L}_n])$ and $d_{nk}^n = \bar{w}^n([L_n, \bar{L}_k])$. Then*

$$(5.6) \quad |(c_{kn}^n L_n U_\ell^{kK}, U_\ell^{nK})| \leq C(\sigma L^A(U) + \sigma^{-1} \|U\|^2),$$

$$(5.7) \quad |(d_{nk}^n \bar{L}_n U_\ell^{nK}, U_\ell^{kK})| \leq C(\sigma L^A(U) + \sigma^{-1} \|U\|^2).$$

Proof. Note that $U_l^{kK} = 0$ on M_0 and $U_l^{nK} = 0$ on M_σ . From (4.29) and (4.30), it follows that

$$(5.8) \quad |c_{kn}^n|_0, |d_{nk}^n|_0 \leq C\sigma^{2m}, \quad 1 \leq k \leq n - 1.$$

Let χ be a C^∞ function defined on $\widetilde{W}(x_0) \ni W(x_0)$ such that $0 \leq \chi \leq 1$, $\chi = 0$ near M_0 , $\chi = 1$ near M_σ and satisfies

$$(5.9) \quad |\chi|_{s, \widetilde{W}(x_0)} \leq C\sigma^{-2ms}, \quad s = 1, 2, \dots$$

Let us write

$$(5.10) \quad (c_{kn}^n L_n U_l^{kK}, U_l^{nK}) = (c_{kn}^n \chi L_n U_l^{kK}, U_l^{nK}) + (c_{kn}^n (1 - \chi) L_n U_l^{kK}, U_l^{nK}).$$

By integration by parts we get, from the estimates in (5.8), that

$$(c_{kn}^n \chi L_n U_l^{kK}, U_l^{nK}) = -(U_l^{kK}, c_{kn}^n \chi \bar{L}_n U_l^{nK}) + \mathcal{O}(\|U\|^2) \lesssim \sigma \cdot \sigma^{4m} \|\chi \bar{L}_n U_l^{nK}\|^2 + \sigma^{-1} \|U\|^2.$$

Using the fact that $\chi = 0$ on M_0 , $U_l^{nK} = 0$ on M_σ , and (5.9), we can perform integration by parts for the function $\|\chi \bar{L}_n U_l^{nK}\|^2$ in a standard way. Then we get

$$(5.11) \quad \|\chi \bar{L}_n U_l^{nK}\|^2 = \|\chi L_n U_l^{nK}\|^2 + \sigma^{-4m} \mathcal{O}(\|U_l^{nK}\|^2 + L(U)).$$

Combining (5.10) and (5.11), we get

$$|(c_{kn}^n \chi L_n U_l^{kK}, U_l^{nK})| \lesssim \sigma L^A(U) + \sigma^{-1} \|U\|^2.$$

Similarly, we can estimate $(c_{kn}^n (1 - \chi) L_n U_l^{kK}, U_l^{nK})$. This proves (5.6). The proof of (5.7) is similar \square

For each small $\rho > 0$, we set

$$S_\rho = \{(y', y_n) ; |y_{2n}| < \rho\} \cap \widetilde{W}(x_0),$$

and let $\lambda_\rho = \psi_{\rho\delta_0} \circ D_{x_0}^{-1}(y)$ be the plurisubharmonic weight functions constructed in Theorem 4.5, where $\delta_0 = \varepsilon\phi(x_0)^{2m}$.

Lemma 5.6. *For each k , $1 \leq k \leq n-1$, set $d_{1k}^n = \bar{\omega}^n([L_1, \bar{L}_k])$ and let $x = D_{x_0}^{-1}(y)$. Then on S_ρ we have*

$$(5.12) \quad |d_{11}^n(\bar{L}_n \lambda_\rho)(y)| \lesssim \sigma^{1/2} \tau(x, \delta_0)^2 \tau(x, \rho \delta_0)^{-2},$$

$$(5.13) \quad |d_{1k}^n(\bar{L}_n \lambda_\rho)(y)| \lesssim \sigma^{1/2} \rho^{-1/2} \tau(x, \delta_0) \tau(x, \rho \delta_0)^{-1+\gamma}, \quad 2 \leq k \leq n-1.$$

Proof. Note that $d_{1k}^n = \bar{\omega}_A^n([L_1^A, \bar{L}_k^A])$, where ω_A^n and L_k^A 's are defined in (4.27) and (4.28). Therefore it follows that

$$(5.14) \quad d_{1k}^n(y) = \sigma^{1/2} \omega^n([L_1, \bar{L}_k])(y) + \mathcal{O}(|A|_1).$$

Now assume that $y = (y', y_{2n}) \in S_\rho$. Then $x = (x', t) = D_{x_0}^{-1}(y)$ satisfies

$$|t| \leq \rho \varepsilon \phi(x_0)^{2m} = \rho \delta_0.$$

Let $\mu_k(x)$ be defined as in (4.33). Then, by functoriality, it follows that

$$\begin{aligned} & \omega^n([L_1, \bar{L}_k])(y) (\bar{L}_n \lambda_\rho)(y) \\ &= \mu_n(x)^{-1} \omega^n([\mu_1(x) H(L_1^\nu), \mu_k(x) \overline{H(\bar{L}_k^\nu)}])(x) \mu_n(x) \bar{L}_n^\nu \psi_{\rho \delta_0}(x) \\ &= \mu_1(x) \mu_k(x) \omega^n([L_1^\nu, \bar{L}_k^\nu])(x) \bar{L}_n^\nu \psi_{\rho \delta_0}(x) + \mathcal{O}(x_{2n} \bar{L}_n^\nu \psi_{\rho \delta_0}(x)), \end{aligned}$$

where we have used the notation in (4.9) and the expression of r_k in (4.10).

Note that $\tau(x, \delta_0) \approx \mu(x, \delta_0) \approx \mu(x, \varepsilon \phi(x)^{2m})$, and $\mu_k(x) \approx \delta_0^{1/2}$ on $\widetilde{W}(x_0)$. Since $x \in S(\rho \delta_0)$ (in x -coordinates), it follows from Proposition 3.8 and the property (iii) of Theorem 3.6 that (set $\delta = \rho \delta_0$ there)

$$(5.15) \quad \begin{aligned} |\mu_1(x)^2 \omega^n([L_1^\nu, \bar{L}_1^\nu])(x) \bar{L}_n^\nu \psi_{\rho \delta_0}(x)| &\lesssim \tau(x, \delta_0)^2 (\rho \delta_0)^2 \tau(x, \rho \delta_0)^{-2} (\rho \delta_0)^{-1} \\ &= \tau(x, \delta_0)^2 \tau(x, \rho \delta_0)^{-2}, \end{aligned}$$

and for $2 \leq k \leq n-1$, we have

$$(5.16) \quad \begin{aligned} & |\mu_1(x) \mu_k(x) \omega^n([L_1^\nu, \bar{L}_k^\nu])(x) \bar{L}_n^\nu \psi_{\rho \delta_0}(x)| \\ &\lesssim \tau(x, \delta_0) \delta_0^{1/2} \cdot (\rho \delta_0)^{1/2} \cdot \tau(x, \rho \delta_0)^{-1+\gamma} (\rho \delta_0)^{-1} \\ &\lesssim \rho^{-1/2} \tau(x, \delta_0) \tau(x, \rho \delta_0)^{-1+\gamma}. \end{aligned}$$

Assuming that A satisfies (4.26), it follows that $|A|_1 \leq \varepsilon_0 \rho^2$ if $y \in S_\rho$. Also, it follows from the property (iii) of Theorem 3.6 that

$$(5.17) \quad |x_{2n} \bar{L}_n^\nu \psi_{\rho \delta_0}(x)| \lesssim 1, \quad x \in S(\rho \delta_0).$$

If we combine (5.14)–(5.17), then (5.12) and (5.13) follow. \square

We now want to prove Theorem 5.1. Assume $U = \sum_{l=1}^n \sum_{|J|=q} U_l^J \bar{\omega}^J \cdot L_l \in \mathcal{B}^q(S_{\varepsilon, \sigma}, A)$ with $\text{supp } U \subset W(x_0)$. Then from (4.40) and (4.41) it follows that

$$T^*U = D'_q U = BU + C|U|,$$

where

$$(5.18) \quad BU = - \sum_{l=1}^n \sum_{|K|=q-1} \sum_{j=1}^n (L_j U_l^{jK}) \bar{\omega}^K \cdot L_l.$$

Also, (2.6) shows that

$$(5.19) \quad SU = D_{q+1}U = AU + C|U|,$$

where

$$(5.20) \quad AU = \sum_{l=1}^n \sum_{|J|=q} \sum_{j=1}^n (\bar{L}_j U_l^J) \bar{\omega}^{jJ} \cdot L_l.$$

Combining (5.18)–(5.20), we see that

$$(5.21) \quad \|AU\|^2 + \|BU\|^2 \leq 2\|SU\|^2 + 2\|T^*U\|^2 + C\|U\|^2.$$

Let us write $U = U' + U''$, where

$$U' = \sum_{l=1}^n \sum_{\substack{|J|=q \\ n \in J}} U_l^J \bar{\omega}^J \cdot L_l, \quad U'' = \sum_{l=1}^n \sum_{\substack{|J|=q \\ n \notin J}} U_l^J \bar{\omega}^J \cdot L_l,$$

and set

$$\begin{aligned} L(U') &= \sum_{l=1}^n \sum_{\substack{k=1 \\ n \in J}}^{n-1} \sum_{|J|=q} \|\bar{L}_k U_l^J\|^2 + \sum_{l=1}^n \sum_{\substack{|J|=q \\ n \in J}} \|L_n U_l^J\|^2, \\ L(U'') &= \sum_{l=1}^n \sum_{\substack{k=1 \\ n \notin J}}^{n-1} \sum_{|J|=q} (\|\bar{L}_k U_l^J\|^2 + \|L_k U_l^J\|^2) + \sum_{l=1}^n \sum_{\substack{|J|=q \\ n \notin J}} \|\bar{L}_n U_l^J\|^2. \end{aligned}$$

Then we can write

$$(5.22) \quad \|AU\|^2 + \|BU\|^2 = \|AU''\|^2 + \|BU''\|^2 + \|AU'\|^2 + \|BU'\|^2 + E(U', U''),$$

where $E(U', U'')$ denotes the (sum of) inner products (AU', AU'') and (BU', BU'') .

Note that the Levi-form of M_σ has at least $(n-2)$ -positive eigenvalues and $U'' = 0$ along M_0 . Therefore we may proceed in the standard way as in [10,11] for U'' and we get

$$(5.23) \quad \|AU''\|^2 + \|BU''\|^2 \geq c \left(L(U'') + \int_{M_\sigma} |U''|^2 dS \right).$$

A typical term of $E(U', U'')$ looks like

$$(L_k U_l^{kK}, L_n U_l^{nK}) - (\bar{L}_n U_l^{kK}, \bar{L}_k U_l^{nK}),$$

where $n \notin K$ and $1 \leq k \leq n-1$. Since $U_l^{kK} = 0$ on M_0 and $U_l^{nK} = 0$ on M_σ , we can perform the integration by parts and we get

$$(5.24) \quad (L_k U_l^{kK}, L_n U_l^{nK}) = (\bar{L}_n U_l^{kK}, \bar{L}_k U_l^{nK}) - (e_k \bar{L}_n U_l^{kK}, U_l^{nK}) \\ + (\bar{e}_n L_k U_l^{kK}, U_l^{nK}) + ([L_k, \bar{L}_n] U_l^{kK}, U_l^{nK}).$$

By integration by parts if necessary, the second and the third term of the right side of (5.24) are bounded by $\sigma L(U') + C\sigma^{-1}\|U\|^2$. If we write

$$(5.25) \quad ([L_k, \bar{L}_n] U_l^{kK}, U_l^{nK}) = \sum_{i=1}^n (c_{kn}^i L_i U_l^{kK}, U_l^{nK}) + \sum_{i=1}^n (d_{kn}^i \bar{L}_i U_l^{kK}, U_l^{nK}),$$

where $c_{kn}^i = \omega^i([L_k, \bar{L}_n])$ and $d_{kn}^i = \bar{\omega}^i([L_k, \bar{L}_j])$, then when $i < n$, the each term of right is dominated by $\sigma L(U') + \sigma^{-1}\|U\|^2$, by applying integration by parts if necessary. If $i = n$, then $(d_{kn}^n \bar{L}_n U_l^{kK}, U_l^{nK})$ is bounded by $\sigma L(U') + \sigma^{-1}\|U\|^2$. The remaining term $(c_{kn}^n L_n U_l^{kK}, U_l^{nK})$ can be handled by using Lemma 5.5. Therefore we conclude that

$$(5.26) \quad |E(U', U'')| \leq C(\sigma L(U') + \sigma^{-1}\|U\|^2).$$

Let $\lambda \in C^\infty(\widetilde{W}(x_0))$ with $|\lambda| \leq 1$, and for $f \in C^\infty(W(x_0))$, we define

$$\|f\|_\lambda^2 = \int_{W(x_0)} |f|^2 e^{-\lambda} dV.$$

Combining (5.22)–(5.26), we conclude that

$$(5.27) \quad \|AU\|^2 + \|BU\|^2 \geq c_1 \left(L(U'') + \int_{M_\sigma} |U''|^2 dV \right) \\ + \frac{1}{3} (\|AU'\|_\lambda^2 + \|BU'\|_\lambda^2) \\ - C\sigma L(U') - C\sigma^{-1}\|U\|^2,$$

because $e^{-\lambda} \geq 1/3$.

Now let us estimate $\|AU'\|_\lambda^2 + \|BU'\|_\lambda^2$. As in (4.2.3) of [11], we get

$$(5.28) \quad \|AU'\|_\lambda^2 + \|BU'\|_\lambda^2 = L_\lambda(U') + \sum_{l=1}^n \sum_{j,k=1}^{n-1} \left[(L_k U_l^{kn}, L_j U_l^{jn})_\lambda - (\bar{L}_j U_l^{kn}, \bar{L}_k U_l^{jn})_\lambda \right],$$

where $L_\lambda(U') = \sum_{l=1}^n \sum_{\substack{|J|=q \\ n \in J}} \sum_{k=1}^{n-1} \|\bar{L}_k U_l^J\|_\lambda^2 + \sum_{l=1}^n \sum_{\substack{|J|=q \\ n \in J}} \|L_n U_l^J\|_\lambda^2$.

With the notation $\delta_k = e^\lambda L_k e^{-\lambda}$, $1 \leq k \leq n-1$, we have

$$(5.29) \quad (\delta_k U_l^{kn}, \delta_j U_l^{jn})_\lambda = (L_k U_l^{kn}, L_j U_l^{jn})_\lambda + \mathcal{O} \left(\sigma^{-1/4} \|(L\lambda)U'\|^2 + \sigma^{1/4} \sum_{k=1}^{n-1} \|L_k U_l^{kn}\|^2 \right),$$

where $\|(L\lambda)U'\|^2 = \|\sum_{k=1}^{n-1}(L_k\lambda)U_l^{kn}\|^2$. By the standard integration by parts method, we obtain that

$$(5.30) \quad \|L_k U_l^{kn}\|^2 = \|\bar{L}_k U_l^{kn}\|^2 - \int_{M_0} d_{kk}^n |U_l^{kn}|^2 dS + \mathcal{O}(\sigma L(U') + \sigma^{-1} \|U'\|^2).$$

As in (5.24), we can write, for $1 \leq j, k \leq n-1$, as

$$(5.31) \quad (\delta_k U_l^{kn}, \delta_j U_l^{jn})_\lambda = (\bar{L}_j U_l^{kn}, \bar{L}_k U_l^{jn})_\lambda - (e_k \bar{L}_j U_l^{kn}, U_l^{jn})_\lambda \\ + (\bar{e}_j \delta_k U_l^{kn}, U_l^{jn})_\lambda + ([\delta_k, \bar{L}_j] U_l^{kn}, U_l^{jn})_\lambda.$$

By integration by parts, if necessary, the second and third terms of (5.31) are bounded by $C\sigma L(U') + C\sigma^{-1} \|U'\|^2$. To estimate $([\delta_k, \bar{L}_j] U_l^{kn}, U_l^{jn})_\lambda$, we write

$$(5.32) \quad ([\delta_k, \bar{L}_j] U_l^{kn}, U_l^{jn})_\lambda = \sum_{i=1}^n (c_{kj}^i L_i U_l^{kn}, U_l^{jn})_\lambda + \sum_{i=1}^n (d_{kj}^i \bar{L}_i U_l^{kn}, U_l^{jn})_\lambda \\ + (\bar{L}_j (L_k \lambda) U_l^{kn}, U_l^{jn})_\lambda,$$

where $c_{kj}^i = \omega^i([L_k, \bar{L}_j])$, and $d_{kj}^i = \bar{\omega}^i([L_k, \bar{L}_j])$.

If $i < n$, then $(d_{kj}^i \bar{L}_i U_l^{kn}, U_l^{jn})_\lambda$ is bounded by $\sigma L(U') + C\sigma^{-1} \|U'\|^2$. By integration by parts, we can write

$$(5.33) \quad (c_{kj}^i L_i U_l^{kn}, U_l^{jn})_\lambda = (c_{kj}^i (L_i \lambda) U_l^{kn}, U_l^{jn})_\lambda + \mathcal{O}(\sigma L(U') + \sigma^{-1} \|U'\|^2).$$

If $i = n$, then $(c_{kj}^n L_n U_l^{kn}, U_l^{jn})_\lambda$ is bounded by $\sigma L(U') + C\sigma^{-1} \|U'\|^2$. By integration by parts again, we can write

$$(5.34) \quad (d_{kj}^n \bar{L}_n U_l^{kn}, U_l^{jn})_\lambda = (d_{kj}^n (\bar{L}_n \lambda) U_l^{kn}, U_l^{jn})_\lambda - \int_{M_0} d_{kj}^n U_l^{kn} \bar{U}_l^{jn} e^{-\lambda} dS \\ + \mathcal{O}(\sigma L(U') + \sigma^{-1} \|U'\|^2).$$

Combining (5.28)–(5.34), we obtain that

$$(5.35) \quad \|AU'\|_\lambda^2 + \|BU'\|_\lambda^2 \geq \sum_{j,k=1}^{n-1} \left[(\bar{L}_j (L_k \lambda) U_l^{kn}, U_l^{jn})_\lambda + \sum_{i=1}^{n-1} (c_{kj}^i (L_i \lambda) U_l^{kn}, U_l^{jn})_\lambda \right] \\ + \sum_{j,k=1}^{n-1} (d_{kj}^n (\bar{L}_n \lambda) U_l^{kn}, U_l^{jn})_\lambda \\ - \sum_{j,k=1}^{n-1} \int_{M_0} (d_{kj}^n U_l^{kn} \bar{U}_l^{jn} + C_1 \sigma^{1/4} d_{kk}^n |U_l^{kn}|^2) e^{-\lambda} dS \\ + cL(U') - C_1 \sigma^{-1/4} \|(L\lambda)U'\|^2 - C\sigma^{-1} \|U'\|^2.$$

With the notation

$$\partial\bar{\partial}\lambda = \sum_{j,k=1}^{n-1} \lambda_{jk} \omega^j \wedge \bar{\omega}^k,$$

the first sum of the right hand side of (5.35) is equal to $\sum_{j,k=1}^{n-1} (\lambda_{jk} U_l^{jn}, U_l^{kn})_\lambda$. Now suppose that $|\lambda| \leq 1$ on $\widetilde{W}(x_0)$. Let $\chi(t)$ denote the function $\sigma^{1/4}/3C_1 e^t$ and set $\phi = \chi(\lambda)$. Then $\chi''(t) \geq C_1 \sigma^{-1/4} \chi'(t)^2$, and we get

$$(5.36) \quad \begin{aligned} \sum_{j,k=1}^{n-1} \phi_{jk} t_j \bar{t}_k &= \chi'(t) \sum_{j,k=1}^{n-1} \lambda_{jk} t_j \bar{t}_k + \chi''(t) \left| \sum_{k=1}^{n-1} (L_k \lambda) t_k \right|^2 \\ &\geq \frac{\sigma^{1/4}}{9C_1} \sum_{j,k=1}^{n-1} \lambda_{jk} t_j \bar{t}_k + C_1 \sigma^{-1/4} \chi'(t)^2 \left| \sum_{k=1}^{n-1} (L_k \lambda) t_k \right|^2. \end{aligned}$$

Since M_0 is pseudoconvex and $d_{kk}^n > d > 0$ for $2 \leq k \leq n-1$, it follows that

$$(5.37) \quad - \sum_{j,k=1}^{n-1} \int_{M_0} (d_{kj}^n U_l^{kn} \bar{U}_l^{jn} + C_1 \sigma^{1/4} d_{kk}^n |U_l^{kn}|^2) e^{-\lambda} dS \geq 0,$$

provided that σ is sufficiently small. Thus, if we replace λ by ϕ , then we conclude from (5.36) and (5.37) that

$$(5.38) \quad \begin{aligned} \|AU'\|_\phi^2 + \|BU'\|_\phi^2 &\geq \frac{\sigma^{1/4}}{9C_1} \sum_{j,k=1}^{n-1} (\lambda_{jk} U_l^{jn}, U_l^{kn})_\phi \\ &\quad + \sum_{j,k=1}^{n-1} (d_{kj}^n (\bar{L}_n \lambda) U_l^{kn}, U_l^{jn})_\phi \\ &\quad + cL(U') - C\sigma^{-1} \|U'\|^2. \end{aligned}$$

Now we take the family $\{\lambda_\rho\}_{\rho>0}$ of plurisubharmonic functions with maximal Hessian constructed in Theorem 4.5, and replace λ in (5.38) by these functions. By Lemma 5.6 and the fact that $\tau(x, \delta) \approx \mu(x, \delta)$, we have, for $y = D_{x_0}(x) \in S_\rho$, that

$$(5.39) \quad |(d_{11}^n (\bar{L}_n \lambda_\rho)(y))| \lesssim \sigma^{1/2} \mu(x, \delta_0)^2 \mu(x, \rho \delta_0)^{-2},$$

$$(5.40) \quad |(d_{1k}^n (\bar{L}_n \lambda_\rho)(y))| \lesssim \sigma^{1/2} \rho^{-1/2} \mu(x, \delta_0) \mu(x, \rho \delta_0)^{-1+\gamma}, \quad 2 \leq k \leq n-1.$$

Also, it follows from (4.27) and (4.28) that

$$(5.41) \quad |(d_{j,k}^n (\bar{L}_n \lambda_\rho)(y))| \lesssim \sigma^{1/2} \rho^{-1}, \quad 2 \leq j, k \leq n-1.$$

By virtue of Theorem 4.5 it follows that (assuming that we first take $\rho_0 \leq \sigma^{2m}$ for sufficiently small σ) there is, for each $0 < \rho \leq \sigma^{2m}$, $\lambda = \lambda_\rho$ such that

$$(5.42) \quad \sum_{j,k=1}^{n-1} \lambda_{j,k}(y) U_l^{jn}(y) \bar{U}_l^{jn}(y) \approx \sum_{k=1}^{n-1} |U^{kn}(y)|^2 \mu_k^2(x) (\mu_k^\rho(x))^{-2}, \quad y \in S_\rho,$$

where $\mu_k(x)$ and $\mu_k^\rho(x)$ are defined in (4.33) and (4.34), respectively.

Combining (5.38)–(5.42), we then have that

$$\begin{aligned} \|AU'\|^2 + \|BU'\|^2 &\geq c\sigma^{1/2} \sum_{j,k=1}^{n-1} (\lambda_{jk} U_i^{jn}, U_i^{kn}) + cL(U') \\ &\geq c(\sigma^{1/2} \rho^{-2/m} \|U'\|^2 + L(U')), \end{aligned}$$

because $\mu(x, \delta_0)^2 \mu(x, \rho\delta_0)^{-2} \gtrsim \rho^{-1/m}$ by (3.34). Then, by the theorem of Catlin [2], the subelliptic estimates of order $1/m$ holds for U' and hence we get

$$(5.43) \quad \sigma^{1/2} \| \|U'\| \|_{1/m}^2 + L(U') \leq C(\|AU'\|^2 + \|BU'\|^2) + \|U'\|^2.$$

Combining (5.21), (5.27) and (5.43), we conclude that

$$(5.44) \quad \|SU\|^2 + \|T^*U\|^2 \geq c \left(\sigma^{1/2} \| \|U'\| \|_{1/m}^2 + L(U) + \int_{M_\sigma} |U''|^2 dS \right) - C\sigma^{-1} \|U\|^2.$$

If $n = 2$, then $U'' = 0$ on M_σ , and if $n \geq 3$, then we have at least one positive eigenvalue. In this case, we apply Lemma 5.3 for f replaced by U'' and get

$$(5.45) \quad \sigma^{1/2} \| \|U''\| \|_{1/2}^2 \leq CL(U'').$$

Combining (5.4), (5.44) and (5.45), we conclude that

$$\sigma^{-3} \|U\|^2 + L(U) + \sigma^{1/2} \| \|U\| \|_{1/m}^2 \leq C(\|SU\|^2 + \|T^*U\|^2),$$

for all $U \in \mathcal{B}^q(S_{\varepsilon,\sigma}; A)$, provided σ is sufficiently small.

For the estimates of the non-tangential derivatives of U , we note that $L_n^A = \partial/\partial y_{2n} + X$, where $X = \sum_{j=1}^{2n-1} b_j(y) \partial/\partial y_j$. Therefore a standard argument yields the inequality

$$(5.46) \quad \left\| \left\| \frac{\partial f}{\partial y_{2n}} \right\| \right\|_{-1+1/m}^2 \leq C \left(1 + \sum_{j=1}^{2n-1} |b_j|_{\widetilde{W}(x_0), 2n+3}^2 \right) (\| \|f\| \|_{1/m}^2 + \| \bar{L}_n f \|^2 + \|f\|^2)$$

for all $f \in C_0^\infty(\widetilde{W}(x_0))$. This inequality can be applied with $f = U_l^J$ to obtain (5.1) from (5.46). This completes the proof of Theorem 5.1. \square

We now define Sobolev spaces for sections of $\Gamma^{0,q}(S_{\varepsilon,\sigma}; A)$. Recall that the open sets $B_b(x_0)$ satisfy (4.3) and (4.4) for each $x_0 \in M$. Choose a set $T_\sigma = \{x_i^\sigma \in M; i \in I\}$ such that the sets $B_{c\sigma/2}(x_i^\sigma)$, $i \in I$, cover $S_{\varepsilon,\sigma}$, and such that no two points x_i^σ and x_j^σ satisfy $|x_i^\sigma - x_j^\sigma| \leq c\sigma/4$, where $|\cdot|$ is the distance function on $S_{\varepsilon,\sigma}$. It follows that the sets $W(x_i^\sigma)$, $i \in I$, cover $S_{\varepsilon,\sigma}$ and that there exists an integer \tilde{N} such that no point of $S_{\varepsilon,\sigma}$ lies in more than \tilde{N} open sets $W(x_i^\sigma)$. Furthermore, there exist functions ζ_i, ζ'_i (that are independent of y_{2n}) $\in C_0^\infty(W(x_i^\sigma))$ such that $\sum_{i \in I} \zeta_i^2 \equiv 1$, and if $x \in \text{supp } \zeta_i$, then

$$(5.47) \quad \zeta'_i \equiv 1 \text{ in } B_{c'\sigma}(x),$$

and that both ζ_i and ζ'_i satisfy

$$(5.48) \quad |\zeta_i|_{k,W(x_i^\sigma)} + |\zeta'_i|_{k,W(x_i^\sigma)} \leq C_k \sigma^{-k}.$$

Now, let F be any section of $\Gamma^{0,q}(S_{\varepsilon,\sigma}; A)$. We define

$$\|F\|_{k,A}^2 = \sum_{i \in I} \|\zeta_i F\|_{k,A,W(x_i^\sigma)}^2,$$

where

$$\|\zeta_i F\|_{k,A,W(x_i^\sigma)}^2 = \sum_{l=1}^n \sum_{|J|=q} \|\zeta_i F_l^J\|_{k,W(x_i^\sigma)}^2,$$

and $F = \sum_{l=1}^n \sum_{|J|=q} F_l^J \bar{\omega}_A^J \cdot L_j^A$ is the decomposition of F in terms of the frame of $W(x_i^\sigma)$.

Moreover, the Sobolev norm $\|\cdot\|_{k,W(x_i^\sigma)}$ is taken with respect to the y -coordinates of $W(x_i^\sigma)$. We define $H_k^{0,q}(S_{\varepsilon,\sigma}; T_A^{1,0})$ to be the set of all sections F of $\Gamma^{0,q}(S_{\varepsilon,\sigma}; A)$ for which $\|F\|_{k,A} < \infty$. If we define $L_q^2(S_{\varepsilon,\sigma}; T_A^{1,0})$ to be the set of all $F \in \Gamma^{0,q}(S_{\varepsilon,\sigma}; A)$ such that $\|F\|^2 < \infty$, then it is obvious that the norms $\|\cdot\|$ and $\|\cdot\|_{0,A}$ are equivalent on $L_q^2(S_{\varepsilon,\sigma}; T_A^{1,0})$. Since $\mathcal{A}(S_{\varepsilon,\sigma}) \subset \Gamma^{0,1}(S_{\varepsilon,\sigma}; 0)$, we define $\|A\|_k = \|A\|_{k,0}$ for $A \in \mathcal{A}(S_{\varepsilon,\sigma})$, and $H_k(S_{\varepsilon,\sigma}; \mathcal{A})$ to be the set of $A \in \mathcal{A}(S_{\varepsilon,\sigma})$ such that $\|A\|_k < \infty$.

We want to get an estimate in global form. Define $Q(U, U) = \|T^*U\|^2 + \|SU\|^2$. By using the partition of unity as defined above satisfying (5.47) and (5.48), and from the estimates in Theorem 5.1, we obtain the following

Corollary 5.7. *There exist a fixed small σ and a constant $\varepsilon_1 > 0$ such that for all $\varepsilon, 0 < \varepsilon < \varepsilon_1$ and all $U \in \text{Dom}(T^*) \cap \text{Dom}(S)$,*

$$(5.49) \quad \|U\|^2 \leq CQ(U, U).$$

Now let us fix $\sigma > 0$ satisfying Corollary 5.7, and set $W(x_0) = W_\sigma(x_0)$. Using Theorem 5.1 and the standard ‘‘bootstrap’’ method, we can get regularity estimates for the linearized equation. The proof is similar to that in Section 9 of [5]. Here we use $1/m$ subelliptic estimates instead of $1/2$ subelliptic estimates. Set $\square = D_q D_q^* + D_{q+1}^* D_{q+1}$. In the sequel, we assume that A satisfies

$$(5.50) \quad \|A\|_{2n+3} \leq \varepsilon_0.$$

Theorem 5.8. *Suppose that U is the solution of $\square U = G$, where $G \in H_k^{0,q}(S_\varepsilon; T_A^{1,0})$ for all $k > 0$. Then*

$$\|D_q^* U\|_k + \|D_{q+1} U\|_k \lesssim \|G\|_k + (1 + \|A\|_{k+2}) \|G\|_{n+2}.$$

Now set $E = D_{q+1}^* D_{q+1} U$. Then we have the following estimates for the error term E [5, Theorem 10.3].

Theorem 5.9. *Suppose that $\square U = G$, where $D_{q+1} G = 0$ and $G \in H_k^{0,q}(S_\varepsilon; T_A^{1,0})$ for all k . Then $E = D_{q+1}^* D_{q+1} U$ satisfies*

$$(5.51) \quad \|E\|_{k-1} \lesssim \|G\|_k \|F^A\|_{n+1} + \|G\|_{n+1} \|F^A\|_k \\ + (1 + \|A\|_{k+2}) \|G\|_{n+1} \|F^A\|_{n+1}.$$

Note that F^A is D_3 -closed. Since $q = 2$, we immediately obtain

Corollary 5.10. *If U is the solution with respect to \mathcal{L}^A of $\square U = F^A$, then $V' = D_2^* U$ satisfies for all $k = n+1, n+2, \dots$*

$$(5.52) \quad \|D_2 V' - F^A\|_k \lesssim \|F^A\|_{k+1} \|F^A\|_{n+1} + \|A\|_{k+3} \|F^A\|_{n+1}^2.$$

6. EXTENSION OF CR STRUCTURES

In this section we will prove Theorem 1.1 and Theorem 1.4, using the estimates in Section 5. First, we describe the nonlinear extension operator. For details, one can refer to Section 11 of [5].

If $A \in \mathcal{A}(S_{\varepsilon, \sigma})$ is sufficiently small and if we set $P_A(\bar{L}) = \bar{L} + A(\bar{L})$, then $\bar{\mathcal{L}}_A = \{P_A(\bar{L}) ; \bar{L} \in \bar{\mathcal{L}}\}$. If we set $Q_A(\omega) = \omega - A^* \omega$, then $\Lambda_A^{1,0} = \{Q_A(\omega) ; \omega \in \Lambda^{1,0}(\mathcal{L})\}$. We define a nonlinear operator $\Phi : \mathcal{A}(S_{\varepsilon, \sigma}) \rightarrow \Gamma^{0,2}(S_{\varepsilon, \sigma})$ by

$$(6.1) \quad \Phi(A)(\bar{L}', \bar{L}'', \omega) = Q_A(\omega)([P_A(\bar{L}'), P_A(\bar{L}'')]).$$

Obviously, if $\Phi(A) = 0$, then $\bar{\mathcal{L}}_A$ is an integrable almost complex structure on $S_{\varepsilon, \sigma}$.

Note that there is a natural map $\mathcal{P}_A : \Gamma_A^{0,2} \rightarrow \Gamma^{0,2}$ defined by

$$(\mathcal{P}_A B)(\bar{L}_1, \bar{L}_2, \omega) = B(P_A(\bar{L}_1), P_A(\bar{L}_2), Q_A(\omega)), \quad B \in \Gamma_A^{0,2}.$$

Therefore it follows from the definition of F^A in (2.5) that $\Phi(A) = \mathcal{P}_A(F^A)$. We also note that if d and A are small sections of \mathcal{A} on $S_{\varepsilon, \sigma}$, then there exist sections $\Delta_{A,d}^+$ and $\Delta_{A,d}^-$ of $\Lambda_A^{0,1} \otimes T_A^{1,0}$ and $\Lambda_A^{0,1} \otimes T_A^{0,1}$, respectively, such that

$$P_{A+d}(\bar{L}) = P_A(\bar{L}) + \Delta_{A,d}^+(P_A(\bar{L})) + \Delta_{A,d}^-(P_A(\bar{L})).$$

Similarly, there exist sections $\delta_{A,\delta}^+$ and $\delta_{A,\delta}^-$ of $\text{Hom}(\Lambda_A^{1,0}, \Lambda_A^{1,0})$ and $\text{Hom}(\Lambda_A^{1,0}, \Lambda_A^{0,1})$, respectively, such that

$$Q_{A+d}(\omega) = Q_A(\omega) - \delta_{A,d}^+(Q_A(\omega)) - \delta_{A,d}^-(Q_A(\omega)).$$

Then it follows that $\Delta_A^\pm(d) = \Delta_{A,d}^\pm$ both depend linearly on d and that the coefficients depend smoothly on A , and the mapping $d \rightarrow \Delta_A(d) = \Delta_A^+(d) + \Delta_A^-(d)$ is invertible. Then $\Phi'(A)(d)$, as an element of $\Gamma^{0,2}$, satisfies

$$(6.2) \quad \Phi'(A)(d) = (\mathcal{P}_A \circ D_2^A \circ \Delta_A^+)(d) - \mathcal{P}_A(h_A(d)(F^A)),$$

where $h_A(d) : T_A^{1,0} \rightarrow T_A^{1,0}$ denotes the adjoint of $\delta_A^+(d) : \Lambda_A^{1,0} \rightarrow T_A^{1,0}$. Since $\Phi(A) = \mathcal{P}_A(F^A)$, we let U_A be the solution of $\square U_A = -F^A$, and then set $V_A = (D_2^A)^* U_A$ and $d_A = \Delta_A^{-1}(V_A)$. Using the error estimates in Theorem 5.9 and Corollary 5.10, we then obtain the following good estimates for the approximate solution of $\Phi(A) + \Phi'(A)(d) = 0$.

Theorem 6.1. *Suppose that $A \in H_k(S_{\varepsilon, \sigma}, \mathcal{A})$ for all k . Then there exists $d_A \in H_k(S_{\varepsilon, \sigma}, \mathcal{A})$ for all k such that if $k \geq n+2$,*

$$(6.3) \quad \|d_A\|_k \lesssim \|\Phi(A)\|_k + \|A\|_{k+2} \|\Phi(A)\|_{n+2},$$

$$(6.4) \quad \|\Phi(A) + \Phi'(A)(d_A)\|_{k-1} \lesssim \|\Phi(A)\|_k \|\Phi(A)\|_{n+2} + \|A\|_{k+2} \|\Phi(A)\|_{n+2}^2.$$

Note that properties (6.3) and (6.4) of the nonlinear operator Φ are the crucial ingredients in the application of simplified Nash-Moser iteration process [15].

We recall that F^A vanishes to infinite order along M_0 (in x -coordinates !). This can be stated in y -coordinates as follows. The proof is similar to that of Lemma 6.2 in [5].

Lemma 6.2. *Suppose that there exists a section $F \in \Gamma^{0,2}(\overline{\Omega}^+)$, where $\overline{\Omega}^+ = \{(x, t) \in \Omega ; 0 \leq t < 1\}$ such that F and all its derivatives vanish to infinite order along M . Then for all $k, N = 0, 1, 2, \dots$, and all $x_0 \in M$,*

$$(6.5) \quad |F^0|_{k,W(x_0)} \leq C_{k,N} \varepsilon^N \varphi(x_0)^N,$$

where F^0 means that F is written out in $W(x_0)$ according to the frames $L_1^0, \dots, L_n^0, \omega_0^1, \dots, \omega_0^n$ of \mathcal{L}^0 .

We can now prove the main theorems of this paper:

Proof of Theorem 1.1. We will show that $\|\Phi(0)\|_D < b$ for the small $b > 0$ and the integer D , which are appeared in the variant of Nash-Moser theorem [15]. The rest properties for the $\Phi(A)$ in the hypothesis of Nash-Moser theorem can be proved using the relations in (6.3) and (6.4), and the estimates for \square operator in Section 5.

Note that (5.48) and (6.5) imply that for each $i \in I$,

$$\|\zeta_i F^0\|_{k,0}^2 \leq C_{k,N} \varepsilon^N \varphi(x_i^\sigma)^N,$$

where ζ_i 's are defined before (5.47). After summing up over x_i^σ , we get

$$(6.6) \quad \|F^0\|_{k,0,\Phi}^2 \leq C_{k,N} \sum_{i \in I} \varphi(x_i^\sigma)^N \varepsilon^N.$$

Since the choice of points that was made before (5.48) shows that the balls $B_{c\sigma/8}(x_i^\sigma)$, $i \in I$, are all disjoint, we can obtain an upper bound on $N(l)$, which is defined to be the number of $i \in I$ such that $2^{-l-1} \leq \varphi(x_i^\sigma) < 2^{-l}$. In fact, in terms of the $\langle \cdot, \cdot \rangle_0$ -metric introduced at the end of Section 2, the volume of $B_{c\sigma/8}(x_i^\sigma)$ is roughly bounded below by $\varepsilon^{n+1} \sigma^{2n-1+2m} \varphi(x_i^\sigma)^{2m(n+2)} \sim \varepsilon^{n+1} \sigma^{2n-1+2m} \cdot 2^{-2lm(n+2)}$, and the $\langle \cdot, \cdot \rangle_0$ -volume of the region in $S_{\varepsilon,\sigma}$ with $2^{-l-1} \leq \varphi(x) \leq 2^{-l}$ is roughly bounded above by $\varepsilon \sigma^{2m} \cdot 2^{-2ml}$. Thus, we conclude that

$$(6.7) \quad N(l) \lesssim \varepsilon^{-n} \sigma^{-(2n-1)} 2^{2ml(n+1)}.$$

Thus (6.6) and (6.7) imply that if $N = 2ml(n+1) + 1$, then

$$\|\Phi(A_0)\|_k = \|F^0\|_{k,0} \lesssim C_k \cdot \varepsilon$$

for sufficiently small ε . In particular, if we set $k = D$ and choose ε to be sufficiently small, then it follow that $\|\Phi(A)\|_D < b$. \square

Proof of Corollary 1.3. Clearly, $\overline{M} \subset bD$ is a CR manifold satisfying all the conditions of Theorem 1.1. Therefore we can extend the given CR structure on \overline{M} to the outside of D by Theorem 1.1. Then, by virtue of Theorem 2.2, the extended CR structure can be patched smoothly with the given complex structure on D . Therefore Corollary 1.3 follows from the Newlander-Nirenberg theorem. \square

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Department of Mathematics
Sogang University
Seoul, 121–742
KOREA

E-mail address : shcho@ sogang.ac.kr