HODGE COHOMOLOGY
OF
NEGATIVELY CURVED MANIFOLDS
by

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#### Abstract

The Hodge Laplacian acting on differential forms is examined for a certain class of complete Riemannian manifolds of dimension $n$. This class consists of the interiors of compact manifolds with boundary, each endowed with a conformally compact' metric. Such a metric, by definition, is of the form $g=\rho^{-2} h$ where $h$ is any smooth nonsingular metric and $\rho$ is a defining function for the boundary. These manifolds are all negatively curved near infinity; examples include the hyperbolic space $\mathbb{H}^{n}$ and those of its quotients which have no cusps.

A parametrix is constructed for the Laplacian acting on the space of $L^{2} k$-forms for all degrees $k \neq \frac{1}{2} n$, $\frac{1}{2}(n+1)$. Thus, in these degrees the Laplacian is Fredholm, and in particular its null space is finite dimensional. The space of $L^{2}$ harmonic $k$-forms is then identified with the relative and absolute deRham cohomology of the manifold when $k<\frac{1}{2}(n-1)$ and $k>\frac{1}{2}(n+1)$, respectively. It is


also shown that the range of the Laplacian is closed when $\mathrm{k}=\mathrm{n} / 2$, although its null space is of infinite dimension.

The parametrix construction is microlocal: the Laplacian should be thought of here as a degenerate elliptic operator on a compact manifold, and a space of pseudodifferential operators large enough to contain its Green operator is defined and studied. A fairly complete calculus, including $L^{2}$ continuity properties, is developed along the way.

Thesis Supervisor: Richard B. Melrose
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Finally, well deserving a paragraph of his own is Richard Melrose. No amount of gratitude on my part would suffice for all that he has done on my behalf.
"Don't brood too much on the superiority of the unseen to the seen. It's true, but to brood on it is mediaeval."

- from Howards End
by E. M. Forster


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## Introduction

The many recent advances in the analysis of partial differential equations notwithstanding, the study of general linear elliptic operators on unbounded domains remains as yet inchoate. The lack of a systematic theory is due no doubt to the fact that the geometry of the domain and the asymptotic degeneracies of the equation each can radically affect the analytic nature of the operator. A natural class of problems in which the geometry and analysis are rather intimately intertwined is in the study of Laplace's operator on a complete Riemannian manifold, acting either on functions or differential forms. Even this case is far from well understood.

A typical question here is whether, on a complete noncompact Riemannian manifold $M$, there exist functions, or k-forms, which are both harmonic and square-summable (or bounded). As was widely suspected from the classical separation of Riemann surfaces into those which do or do not admit bounded harmonic functions, the operative feature is the curvature. (Of course, the $L^{2}$ class is not well defined in the conformal geometry of Riemann surfaces, whereas boundedness is an invariant concept, so long as attention is restricted to functions.) The first general theorem of this type in higher dimensions is due to Yau
[28], and asserts the nonexistence of nonconstant bounded harmonic functions on complete manifolds of nonnegative Ricci curvature.

Negative curvature produces the opposite effect, as may be discerned from the harmonic analysis of the hyperbolic space $\mathbb{H}^{n}$. Although this space has no $L^{2}$ harmonic functions, it admits many bounded harmonic functions--in fact, by a Poisson-type representation theorem, there is one corresponding to each continuous function on $S^{n-1}=$ $\partial \mathbb{H}^{\mathrm{n}}$ (and others with less regular boundary values).

Much attention over the past several years has focused on the Laplacian of complete negatively curved manifolds, and some quite general results have been obtained. Let us discuss some of these, first for the Laplacian on functions. Here the concern is only with bounded harmonic functions since there are no nontrivial $L^{2}$ ones. In fact, if $u \in L^{2} \cap$ dom $\Delta$ is harmonic, an integration by parts-which relies on the completeness of $M$ to ensure the lack of a boundary term-implies that $d u=0$. Thus $u$ is constant, and must vanish if the volume of $M$ is infinite.

If $M$ is simply connected and has all sectional curvatures $K_{M}$ negative, then the theorem of Cartan-Hadamard asserts that $\exp _{p}: T_{p} M \longrightarrow M$ is a diffeomorphism for each $p \in M$. It is possible to define a sphere at infinity, $S_{\infty}$, for such manifolds, cf. Eberlein-O'Neill
[15]. Each point $x \in S_{\infty}$ is identified with a set of geodesics all asymptotic to one another, and these are said to converge to $x . \quad S_{\infty}$ itself is the aggregate of such mutually asymptotic classes of geodesics. There is a natural topology which makes $\bar{M}=M U S_{\infty}$ into a compact topological manifold with boundary (in fact, a disc). The asymptotic Dirichlet problem may then be posed: for any given continuous function $f$ on $S_{\infty}$, find a function $u$ on $M$ which is harmonic and assumes the boundary values $f$ asymptotically. It should be remarked that it is insufficient merely to require

$$
\lim _{t \rightarrow \infty} u(r(t))=f(x)
$$

for any geodesic $r(t)$ converging to $x \in S_{\infty}$; instead, the convergence must occur in the topology of $M$, which is stronger than this 'radial' convergence.

This asymptotic Dirichlet problem was posed by Choi in [7] and partially solved. For various technical reasons one must assume that the sectional curvatures are bounded away from zero: $K_{M} \leq-\mathrm{a}^{2}<0$. Choi proved the result with this hypothesis, but also assuming a convexity condition at infinity-namely that for two distinct points $x, y \in S_{\infty}$ there are disjoint geodesically convex neighbourhoods $U_{x}$, $U_{y}$ with $x \in U_{x}, y \in U_{y}$. This extra condition is unnecessary when $M$ is a surface. Mike Anderson finished
the proof [1] by showing that if there is a two-sided bound on curvature, $-b^{2} \leq K_{M} \leq-a^{2}<0$, then $M$ satisfies the convexity condition. Sullivan independently proved the result [27] using probabilistic methods.

This positive resolution has the consequence that $M$ carries many bounded harmonic functions. In fact, the asymptotic maximum principle of Yau asserts that, in this situation, the infimum and supremum of $u$ are attained on $S_{\infty}$, hence coincide with those of $f$. At this juncture it is reasonable to seek a Poisson-type representation formula. This was accomplished in Anderson-Schoen [3].

Their paper also addresses a quite interesting and relevant question: namely, to what extent may $\bar{M}$ be considered a smooth manifold? Inasmuch as there is a homeomorphism between the unit sphere $S_{p} M$ in the tangent space at any point and $S_{\infty}$, one may form the composite map

$$
\mathrm{T}_{\mathrm{pq}}: \mathrm{S}_{\mathrm{p}} \mathrm{M} \longrightarrow \mathrm{~S}_{\infty} \longrightarrow \mathrm{S}_{\mathrm{q}} \mathrm{M}
$$

and study its regularity. Anderson and Schoen prove that this map is Hölder continuous with exponent $\alpha$, with $\alpha=$ $a / b$ depending only on the curvature bounds. Little else is known one way or the other about this question, though it seems likely that $T_{p q}$ is rarely $C^{\infty}$. We note, however, the result of Fefferman [16] concerning the Bergman metric in the interior of a smoothly bounded strictly
pseudoconvex domain. This metric has asymptotically constant holomorphic sectional curvatures, and he proves that each $T_{p q}$ (which is only defined locally) is $C^{\infty}$; in fact, the map from $S_{p} M$ to the boundary of the domain is $C^{\infty}$.

Finally, on a spectral note, it is also known that the spectrum of the Laplacian (on functions) is contained in the interval $\left[(n-1)^{2} a^{2} / 4, \infty\right)$, cf. McKean [21]. In addition, if $M$ (is simply connected and) has sectional curvatures tending uniformly to $-\infty$, its Laplacian has pure point spectrum, $c f$. Donnelly-Li [13].

Much less is known about the Laplacian acting on differential forms. Of course, when $M$ is compact the Hodge theorem provides topological meaning to the dimension of the space of harmonic forms in each degree. When $M$ is complete but noncompact, it is natural to allow the Laplacian to act on

$$
L^{2} \Omega^{k}=\left\{\omega \in \Omega^{k}=\Gamma\left(\Lambda^{k_{M}}\right): \int|\omega|^{2}<\infty\right\}
$$

the space of $L^{2} k$-forms over $M$. The subspace of harmonic forms is denoted

$$
\#^{k}=\left\{\omega \in L^{2} \Omega^{k} \cap \operatorname{dom} \Delta: \Delta \omega=0\right\}
$$

Recall that $\omega \in \operatorname{dom} \Delta$ means $d \omega \in L^{2} \Omega^{k+1}, \delta \omega \in L^{2} \Omega^{k-1}$,
and $\delta d \omega, d \delta \omega \in L^{2} \Omega^{k}$. The (by now) classical theorem of Andreotti-Vesentini [4] allows one to integrate by parts, as on a compact manifold, to deduce that

$$
\mathscr{H}^{\mathrm{k}}=\left\{\omega \in \mathrm{L}^{2} \Omega^{\mathrm{k}}: \mathrm{d} \omega=\delta \omega=0\right\}
$$

provided $M$ is complete. Notice that this implies $\omega \in$ dom $\Delta$ if $\omega \in L^{2} \Omega^{k}$ and $\Delta \omega=0$ pointwise. We shall call $\mathscr{H}^{k}$ the Hodge cohomology space of degree $k$.

Another similar, but perhaps more widely studied, space is the $L^{2}$ cohomology space of degree $k$ :

$$
\mathrm{L}^{2} \mathrm{H}^{\mathrm{k}}=\left\{\omega \in \mathrm{L}^{2} \Omega^{\mathrm{k}}: \mathrm{d} \omega=0\right\} / \mathrm{d}\left\{\eta \in \mathrm{~L}^{2} \Omega^{\mathrm{k}-1}: \mathrm{d} \eta \in \mathrm{~L}^{2} \Omega^{\mathrm{k}}\right\}
$$

Since each $\omega \in \mathscr{H}^{k}$ is closed, there is a natural map $\mathscr{H}^{k} \longrightarrow L^{2} H^{k}$ for each $k$. However, unlike the compact setting, the two spaces are rarely isomorphic. A gauge of their relationship is afforded by Kodaira's weak Hodge decomposition [9]:

$$
\mathrm{L}^{2} \Omega^{\mathrm{k}}=\overline{\mathrm{d}\left\{\mathrm{~L}^{2} \Omega^{\mathrm{k}-1} \cap \operatorname{dom} \mathrm{~d}\right\}} \oplus \overline{\delta\left\{\mathrm{L}^{2} \Omega^{\mathrm{k}} \cap \operatorname{dom} \delta\right\}} \oplus \not \mathscr{H}^{\mathrm{k}}
$$

which is valid quite generally. The summands here are pairwise orthogonal. From this it is easy to see that

$$
\mathscr{H}^{k} \simeq\left\{\omega \in L^{2} \Omega^{k}: \mathrm{d} \omega=0\right\} / \overline{d\left\{L^{2} \Omega^{k-1} \cap \operatorname{dom} \mathrm{~d}\right\}}
$$

whence the map $\mathscr{H}^{\mathrm{k}} \longrightarrow \mathrm{L}^{2} \mathscr{H}^{\mathrm{k}}$ is always injective.
Furthermore, the two spaces coincide precisely when the range of $d$ is closed.

The $L^{2}$ cohomology will rarely be mentioned again.
The bulk of this dissertation is devoted to the computation of $\mathbb{H}^{k}$ for a certain class of manifolds. It has been suspected for quite a while that if $M$ is simply connected, and its sectional curvatures satisfy $-b^{2} \leq K_{M} \leq-a^{2}<0$, then most of the harmonic spaces $\#^{k}$ should be trivial.

It is straightforward to do the necessary calculations for the hyperbolic space $\mathbb{H}^{n}$, or in fact for any complete rotationally symmetric manifold (which need not satisfy the curvature constraints), see Dodziuk [10], to arrive at the following conclusion: if in polar coordinates the metric of such a manifold has the form $d s^{2}=d r^{2}+f(r)^{2} d \theta^{2}, r \in$ $\mathbb{R}^{+}, \quad \theta \in S^{n-1}$, then

$$
\operatorname{dim} \mathscr{H}^{k}= \begin{cases}1 & \text { if } k=0, n \text { and } \int_{0}^{\infty} f(r)^{n-1} d r<\infty \\ \infty & \text { if } k=n / 2 \text { and } \int_{1}^{\infty} \frac{d r}{f(r)}<\infty \\ 0 & \text { otherwise. }\end{cases}
$$

The first line simply asserts that constant functions are in $L^{2}$ iff the volume of $M$ is finite. The second is
more interesting; it reflects the fact that both the kernel of the coboundary operator and the $L^{2}$ norm are conformal invariants on $\Omega^{n / 2}$. The finiteness of the integral means that $M$ is conformally equivalent to a Euclidean disc of finite radius: introduce $R(r)=\exp \int_{1}^{r} \frac{d t}{f(t)}$ as the new radial variable. Then one only need note that there certainly is an infinite dimensional family of smooth harmonic forms on such a disc.

It is striking that in this proposition, the important geometric feature is the growth of the metric rather than the curvature. This is unexpected since the spaces $\mathscr{H}^{k}$ implicitly depend on the derivatives of the metric, by way of the coefficients of the Laplacian. If $g^{\prime}$ is uniformly equivalent to $g, \quad c_{1} g \leq g{ }^{\prime} \leq c_{2} g$, then its Laplacian has domain different than that for $g$ unless g' is uniformly $C^{2}$ close to $g$. In contradistinction, the spaces $L^{2} H^{k}$ are patently stable under such a $C^{0}$ perturbation of the metric.

Further support for the conjectured vanishing of Hodge cohomology outside of the middle degree(s) was given by the discovery of Donnelly-Xavier [14] that, so long as the curvature is (negative and) tightly enough pinched, $|a / b-1|<$ $\epsilon(n, k)$, then $\mathscr{H}^{k}=0$ for all $k$ such that $|k-n / 2|>$ 1/2. Donnelly-Fefferman [12] extended the method to show that for the Bergman metric of a strictly pseudoconvex
domain in $\mathbb{C}^{n}$, where $\mathscr{H}^{k}$ splits as $\oplus_{\mathrm{p}+\mathrm{q}=\mathrm{k}} \mathscr{H}^{\mathrm{p}, \mathrm{q}}$, then

$$
\operatorname{dim} \not \mathscr{H}^{p, q}=\left\{\begin{array}{ll}
0 & p+q \neq n \\
\infty & p+q=n
\end{array} .\right.
$$

Recall that the Bergman metric has holomorphic sectional curvatures asymptotically constant at the boundary of the domain.

It was thus quite surprising when Mike Anderson [2] constructed counterexamples to the conjecture; there exist simply connected manifolds with $-\mathrm{b}^{2} \leq \mathrm{K}_{\mathrm{M}} \leq-\mathrm{a}^{2}<0$ for which $\#^{k}$ is infinite dimensional for any given k. Furthermore, it is seen that the curvature pinching ratio of Donnelly-Xavier is sharp.

Many questions are thereby raised dealing with how one might recognize the negatively curved manifolds with suitable vanishing of the harmonic spaces. It appears that the key to this problem lies in the (intrinsic) regularity of the sphere at infinity. The characterization of those manifolds for which $S_{\infty}$ is smoothly attached is very likely quite difficult, so first step might be to a priori assume some regularity at infinity. Among the many forms this assumption could take, we pursue the course of considering manifolds with metric differing from that of the interior of a compact Riemannian manifold with boundary by
a singular conformal factor; these we call conformally compact.

More explicitly, let $M$ be a compact manifold with boundary, $h$ a metric nondegenerate and smooth up to $\partial M$, and $\rho$ a function defining $\partial M$ (so that $\left.d \rho\right|_{\partial M} \neq 0$ ). Consider the metric $g=\rho^{-2} h$. The singular factor $\rho^{-2}$ has the effect of pushing $\partial M$ to infinity. It is not hard to see that, like the Poincare metric on the unit ball after which it is patterned, this metric is complete and, more importantly, negatively curved near $\partial M$. There are other ways to produce metrics with inverse-square asymptotics, and the correct geometry. In particular, the Bergman metric is obtained as the complex Hessian of

$$
c \log \left(\phi_{1} \rho^{-n-1}+\phi_{2} \log \rho\right)
$$

with $\phi_{1}, \phi_{2}$ smooth, $\phi_{1}$ nonvanishing, and $\rho$ the distance to the boundary. (For the unit ball, $\phi_{1}=1$ and $\left.\phi_{2}=0.\right)$ But, were it to be written as $\rho^{-2} h$, the metric $h$ then partially degenerates at $\partial M$-in fact, precisely in the directions of the $C R$ bundle in $T \partial M$.

Our principal result is the

Theorem (4.8): For the metric $g=\rho^{-2} h$ on $M$, as described above, there are natural isomorphisms

$$
\begin{array}{ll}
\mathscr{H}^{k} \simeq H^{k}(M, \partial M) & k<\frac{n-1}{2} \\
\mathscr{H}^{k} \simeq H^{k}(M) & k>\frac{n+1}{2}
\end{array}
$$

obtained by sending a (necessarily closed) form $\omega \in \mathscr{H}^{k}$ to its de Rham class. By virtue of conformal invariance, $\mathscr{H}^{k}$ is infinite dimensional when $k=n / 2$.

This is proved by constructing a parametrix for $\Delta_{g}$, i.e. a pseudodifferential operator $E$ such that

$$
\Delta_{\mathrm{g}} \mathrm{E}=\mathrm{I}-\mathrm{Q}, \quad \mathrm{E}^{*} \Delta_{\mathrm{g}}=\mathrm{I}-\mathrm{Q}^{*}
$$

where $Q$ and its adjoint $Q^{*}$ are compact and smoothing on $L^{2}(\mathrm{dg})$. From a fairly explicit knowledge of the Schwartz kernels of $Q$ and $Q^{*}$, it is possible to deduce the boundary asymptotics of elements of $\boldsymbol{H}^{k}$. This, in turn, allows the isomorphisms of the theorem to be proved. A closely related 'Poisson' operator solves the asymptotic Dirichlet problem, but due to constraints of time and space this will not be developed here.

The main difficulty in the construction of $E$ arises from the fact that the principal symbol of $\Delta_{g}$ degenerates quadratically at $\partial M$. The Schwartz kernel of $E$ must therefore have a fairly complicated singularity at the corner of $M \times M$. The systematic introduction of singular
coordinates at this corner allows one to adequately 'resolve' this singularity. In effect, the degeneracies are transferred from the Schwartz kernel of the operator in question to the geometry of the new 'blown up' manifold which replaces $M \times M$. This type of construction was first developed by Melrose [22] to deal with differential operators degenerating somewhat less thoroughly than the ones considered here. Melrose-Mendoza [23] studies that class of operators further, in particular proving Fredholm properties for the elliptic elements. Both that work and the present one are part of more general construction, to be contained in the, as yet mythical, Melrose-Mendoza [24]-from which inspiration for the following pages comes.

The interest in these methods as presented here lies perhaps in the ease of their applicability to a natural geometric question (albeit in a somewhat artificial setting). It seems likely that similar techniques will prove effective in other related geometric problems. For instance, Mazzeo-Melrose [20] employs essentially the same construction as here to deduce the meromorphic extension of the resolvent and Eisenstein series for certain quotients of $\mathbb{H}^{n}$.

This paper is organized as follows. The first chapter examines conformally compact manifolds from a differential geometric point of view. Their curvature is seen to have fairly simple behaviour, and in particular is negative at
infinity. Asymptotics of the geodesics are also studied, and it is shown that in a special case the ideal boundary has an intrinsic regularity structure. In the second chapter, the $V_{0}$ operator calculus is developed, including $L^{2}$ continuity of its elements, and the parametrix construction is outlined. In Chapter 3 , a certain model for $\Delta_{g}$ is analyzed. In these circumstances this model is nothing but the (constant curvature) hyperbolic Laplacian; however we require some rather unusual mapping properties of it. Finally, the hard work now complete, the actual construction and proof of the Main Theorem are contained in the brief last chapter.

In conclusion we recall an old conjecture of Heinz Hopf's, for it has stimulated much of the work directed toward the computation of $\mathscr{H}^{k}$ on negatively curved manifolds. One formulation of the classical uniformization theorem classifies surfaces by which constant curvature metric they admit. Motivated both by this and the higherdimensional Gauss-Bonnet Theorem, Hopf proposed the generalization:

Conjecture: If $M^{2 m}$ is compact and admits a metric with strictly negative sectional curvatures, then $(-1)^{m} \times(M)>0$.

The corresponding query for compact flat manifolds is trivial, and the one for compact manifolds of positive curvature is likewise unresolved. The proposed methods of proof are quite different for the positive and negative cases; in the positive case much work has been devoted toward showing that the Gauss-Bonnet integrand-the Pfaffian-is pointwise positive. This was shown to be a hopeless task by Geroch [17] (who gives a local counterexample) and Bourgouignon-Karcher [6] (who give a global one).

Following a similar pattern of history, Anderson's counterexample [2] likewise dashed hopes that the negative case could be resolved by settling the hopefully less complicated issue of determining the $L^{2}$ harmonic spaces on complete noncompact negatively curved manifolds. In fact, Singer had suggested that combining Atiyah's index theorem for covers [5] with suitable vanishing of the Hodge cohomology would settle Hopf's conjecture. This index theorem equates the alternating sum of the Betti numbers of a compact manifold $M$, which of course is just the Euler characteristic, with the alternating sum of certain $L^{2}$ Betti numbers on the universal cover $\tilde{M}$ :

$$
x(M)=\sum(-1)^{k_{b_{k}}(M)=\sum(-1)^{k_{b}} \Gamma_{k}(\tilde{M}), ~(M)}
$$

$\Gamma=\pi_{1}(M)$. These $b_{k}^{\Gamma}$ are nothing but the (normalized)
dimension of the $\mathscr{H}^{k}$. Thus, since if the curvature is negative $\tilde{M}$ is simply connected, if one could show that $\mathscr{H}^{k}=0$ for $k \neq m=\frac{1}{2} \operatorname{dim} M$ for the negatively curved pullback metric on $\tilde{M}$, it follows that

$$
x(M)=\sum(-1)^{k_{b}}{ }_{\mathrm{k}}^{\Gamma}(\tilde{M})=(-1)^{m_{b}}{ }_{\mathrm{m}}^{\Gamma} .
$$

The $b_{k}^{\Gamma}$ are always nonnegative, and it is not hard to show that $b_{m}^{\Gamma}$ is strictly positive, so the conjecture is proved.

This method of proof does work when the curvature is close to being constant-and in particular when $M$ is hyperbolic-by virtue of Donnelly-Xavier's work. We remark also that Anderson's counterexamples are not universal covers of compact manifolds since their isometry groups have no cocompact subgroups. Thus it may well still be true that such periodic covering manifolds have vanishing Hodge cohomology outside the middle degree, but to prove this would likely be quite difficult. It is, I think, a question that warrants attention.

Chapter 1. Conformally Compact Metrics

## A. Definitions

The interior of a compact manifold with boundary $M$ may be endowed with a complete Riemannian metric, and thus becomes a complete open manifold. The general effect is that $\partial M$ is placed at infinity; geodesics take an infinite time to reach it. The seminal example, and model, for us is the Poincare metric on the unit ball $B^{n}$. In local coordinates near the boundary the components of the metric become arbitrarily large. The rate of this blow-up is essentially fixed by balancing the requirements that the metric be complete and that its sectional curvatures be negative and bounded away from zero near $\partial M$. Explicitly, suppose
$h$ is a nondegenerate smooth metric on the (closed) manifold $M$.
(1.1) $\rho$ is a defining function for $\partial M: \rho \geq 0$, $\rho^{-1}(0)=\partial M$, and $d \rho$ is nonvanishing on $N^{*} \partial M \backslash 0$
and set

$$
\mathbf{g}=\rho^{-\sigma} \mathrm{h}, \quad \sigma>0 .
$$

Then, from computations in sections $B$ and $C$ of this chapter
it will follow that $g$ is complete for $\sigma \geq 2$, whereas its sectional curvatures are bounded above by a negative constant only when $\sigma \leq 2$. In fact the curvature diverges to negative infinity if $\sigma$ is strictly less than 2.

This thesis focuses on a study of the analytic properties of

$$
\begin{equation*}
g=\rho^{-2} h \tag{1.2}
\end{equation*}
$$

which is both complete and negatively curved near $\partial M$, and in particular of its Laplacian. Let us however mention that there are other examples of metrics of this general form which are both complete and negatively curved, but which are not "conformally compact"-i.e. as in (1.2) with $\rho$ and $h$ satisfying the hypotheses of (1.1). For example, the Bergman metric on the unit ball in $\mathbb{C}^{n}$ has the correct geometry (and there is a vanishing theorem for its Hodge cohomology), but in Cartesian coordinates it blows up as $\left(1-|z|^{2}\right)^{-1}$ in directions converging to the complex subbundle of $T \partial B^{n}$ and as $\left(1-|z|^{2}\right)^{-2}$ only in the other two directions. It would be quite interesting to understand the ways $\rho$ and $h$ could degenerate so that $\rho^{-2} h$ still is complete and of bounded geometry.

This chapter is devoted to a discussion of the geometry of (M,g), $g$ as in (1.2). In particular, we examine
the behaviour of diverging geodesics fairly thoroughly, as well as asymptotic properties of the curvature tensor. All computations use $\rho$ and $h$ routinely, but note that, in addition, $g=(\phi \rho)^{-2} \phi^{2} h$ for any strictly positive $\phi \in$ $C^{\infty}(M)$. All formulae for the geometric quantities of $g$ must therefore be invariant under such a transformation.

## B. Geodesics at Infinity

Introduce coordinates $\left(z^{1}, \cdots, z^{n}\right)$ near a point $p \in$ $\partial M$ such that $z^{n}=0$ on the boundary and $\partial / \partial z^{n}$ is the inward pointing $h$-unit normal. $z^{\prime}=\left(z^{1}, \cdots, z^{n-1}\right)$ are then coordinates on $\partial M$. It is more convenient to work with the co-geodesic flow, the equations for which, in terms of $z$ and the dual coordinates $\xi_{1}, \cdots, \xi_{n}$ are

$$
\begin{gathered}
\dot{z}^{i}=\rho^{2} h^{i j} \xi_{j} \\
\dot{\xi}_{i}=-\rho \frac{\partial \rho}{\partial z^{i}} h^{p q} \xi_{p} \xi_{q}-\frac{1}{2} \rho^{2} \frac{\partial h^{p q}}{\partial z^{i}} \xi_{p} \xi_{q} \\
=-\frac{1}{\rho} \frac{\partial \rho}{\partial z^{i}}-\frac{1}{2} \rho^{2} \frac{\partial h^{p q}}{\partial z^{i}} \xi_{p} \xi_{q} .
\end{gathered}
$$

Here all indices vary between 1 and $n$, the summation convention on repeated indices is used, and attention is restricted to the energy surface

$$
\begin{equation*}
\rho^{2} h^{i j} \xi_{i} \xi_{j}=1 \tag{1.4}
\end{equation*}
$$

Notice that (1.4) implies

$$
\begin{equation*}
\mathrm{C}_{1} / \rho \leq|\xi| \leq \mathrm{C}_{2} / \rho, \quad|\xi|^{2}=\xi_{1}^{2}+\cdots+\xi_{\mathrm{n}}^{2} \tag{1.5}
\end{equation*}
$$

$C_{1}, C_{2}$ depend only on the largest and largest eigenvalues of the matrix $h^{i j}$.

Now consider a maximally extended geodesic $r(t)=$ $\left(z^{1}(t), \cdots, z^{n}(t)\right)$ and suppose

$$
\rho(r(0))<\epsilon, \quad \dot{\mathbf{z}}^{\mathrm{n}}(0)<0
$$

where $\epsilon$ is sufficiently small. $\gamma$ is, of course, the projection of a Hamiltonian curve ( $z^{1}, \cdots, z^{n}, \xi_{1}, \cdots, \xi_{n}$ ) which solves (1.3). By virtue of the choice of coordinates,

$$
\partial \rho / \partial z^{i}=a_{i}(z) \rho(z), \quad i=1, \cdots, n-1, \quad \partial \rho / \partial z^{n}>0
$$

for some $C^{\infty}$ functions $a_{i}$, and in particular $\rho^{-1} \partial \rho / \partial z^{i}$ is smooth up to $\partial M$. Using (1.5) in (1.3) one concludes

$$
\dot{\xi}_{i}=0(1) \quad i=1, \cdots, n-1
$$

(1.6)

$$
-\mathrm{k}_{1} / \rho \leq \dot{\xi}_{\mathrm{n}} \leq-\mathrm{k}_{2} / \rho
$$

where $k_{1}, k_{2}$ depend only on $\epsilon$. We first prove
(1.7) Lemma: If $r(t)$ is a geodesic as described above, and $\epsilon>0$ is sufficiently small, then $\dot{z}^{n}(t)<0$ for all $t>0$.

Proof: By (1.3) $\dot{z}^{n}=\rho^{2} h^{n i} \xi_{i}=\rho^{2}\left(h^{n \alpha} \xi_{\alpha}+h^{n n} \xi_{n}\right)$ where the sum in $\alpha$ is from 1 to $n-1$. Set $F(t)=h^{n \alpha} \xi_{\alpha}+$ $h^{n n} \xi_{n}$; then $F(0)<0$ by assumption, and we compute

$$
F^{\prime}(t)=h^{n \alpha} \dot{\xi}_{\alpha}+h^{n n} \dot{\xi}_{n}+\frac{\partial h^{n \alpha}}{\partial z^{i}} \dot{z}^{i} \xi_{\alpha}+\frac{\partial h^{n n}}{\partial z^{i}} \dot{z}^{i} \xi_{n}
$$

By (1.5) and (1.3), $\dot{\mathbf{z}}^{\mathbf{i}}=O(\rho)$ so $\quad \dot{\mathbf{z}}^{\mathrm{i}} \xi_{\alpha}, \quad \dot{\mathbf{z}}^{\mathbf{i}} \xi_{\mathrm{n}}$ are $\quad \mathrm{O}(1)$. Thus

$$
\mathrm{F}^{\prime}(\mathrm{t})=\mathrm{h}^{\mathrm{nn}} \dot{\xi}_{\mathrm{n}}+0(1) \leq \mathrm{h}^{\mathrm{nn}} \dot{\xi}_{\mathrm{n}}+\mathrm{C}
$$

where $C$ again depends only on the metric but not on $r$. By applying the bound (1.6) and assuming $\epsilon$ is small enough we may ensure that

$$
\dot{\xi}_{n}(t)<-K \quad \text { if } \quad z^{n}(t)<z^{n}(0)
$$

where $K$ is chosen so that $-\frac{1}{2} K+C<-1$. Now set

$$
\mathrm{t}_{0}=\inf \left\{\mathrm{t}>0: \dot{\mathrm{z}}^{\mathrm{n}}(\mathrm{t}) \geq 0\right\}
$$

If $0<t<t_{0}$ then $z^{n}(t)<z^{n}(0)$, so $\dot{\xi}_{n}(t)<-K$. Finally, if $h^{n n}>1 / 2$ in $\{z: \rho(z)<\epsilon\}$, we conclude that

$$
F^{\prime}(t) \leq h^{n n} \dot{\xi}_{n}+C \leq-\frac{1}{2} K+C<-1
$$

for all $0<t<t_{0}$ Hence $F\left(t_{0}\right)<F(0)<0$, ie. $\dot{z}^{\mathrm{n}}\left(\mathrm{t}_{0}\right)<0$, a contradiction.
(1.8) Proposition: $g$ is a complete metric. In particular, if $\quad \gamma(\mathrm{t})$ is a geodesic with $\rho(\gamma(0))<\epsilon$, $\dot{z}^{\mathrm{n}}(0)<0$, and $\epsilon$ is small enough, then $\gamma$ can be extended to infinite length.

Proof: First observe that the function $z^{n}$ of our coordinate system may be defined on a collar neighbourhood of the full boundary. On this neighbourhood there is a product metric

$$
k=\left(z^{n}\right)^{-2}\left(\left.h\right|_{\partial M}+\left(d z^{n}\right)^{2}\right)
$$

Clearly for some $c, g \geq c k$. This implies that

$$
\begin{aligned}
1 & =\left|\gamma^{\prime}(\mathrm{t})\right|_{\mathrm{g}} \geq \mathrm{c}\left|\dot{z}^{\mathrm{n}}\right| / z^{\mathrm{n}}=-\mathrm{c} \dot{\mathbf{z}}^{\mathrm{n}} / \mathrm{z}^{\mathrm{n}} \\
\Rightarrow \mathrm{t} & =\int_{0}^{\mathrm{t}}\left|\gamma^{\prime}(\tau)\right|_{\mathrm{g}}^{\mathrm{d} \tau} \geq \int_{0}^{\mathrm{t}}(-\mathrm{c}) \frac{\dot{\mathbf{z}}^{\mathrm{n}}(\tau)}{z^{\mathrm{n}}(\tau)} \mathrm{d} \tau
\end{aligned}
$$

which by Lemma (1.7) reduces to

$$
t \geq c \int_{z^{n}(t)}^{z^{n}(0)} \frac{d z^{n}}{z^{n}}=c \log \left(z^{n}(0) / z^{n}(t)\right)
$$

In other words, $t$ tends to infinity as $z^{n}(t)$ tends to 0 , and consequently $r$ has infinite length.

Notice that completeness still obtains if $\left(z^{n}\right)^{-2}$ is replaced by $\left(z^{n}\right)^{-\sigma}, \sigma>2$, in $k$ and the rest of the proof is modified accordingly. Now, as might be hoped, it is also true that geodesics approach definite points of the boundary. More is true:
(1.9) Proposition: $\gamma(t)$ tends to a definite point of the boundary as $t \longrightarrow \infty$. Furthermore, taking $z^{n}$ as the new variable, the reparametrized geodesic is tangent to the h-unit normal $\partial / \partial z^{n}$ at the boundary.

Proof: Since $\dot{\xi}_{\mathrm{n}} \longrightarrow-\infty$, we may use $\xi_{\mathrm{n}}$ as a nonsingular parameter, and then

$$
\begin{aligned}
& \partial z^{i} / \partial \xi_{n}=\rho^{2} h^{i j} \xi_{j} /\left(-2 \frac{\rho}{\rho}+O(1)\right) \\
& =\rho^{3} h^{i j} \xi_{j} /\left(-2 \rho_{n}+O(\rho)\right)=O\left(\rho^{2}\right) . \\
& \mathrm{i}=1, \cdots, \mathrm{n}-1, \quad \rho_{\mathrm{n}}=\partial \rho / \partial \mathrm{z}^{\mathrm{n}} \text {. From (1.5) it follows that } \\
& \left|\partial z^{i} / \partial \xi_{n}\right| \leq C \xi_{n}^{-2} . \\
& \text { Hence } r(t) \text { tends to }\left(z_{\infty}^{1}, \cdots, z_{\infty}^{n-1}, 0\right) \text {, where } \\
& z_{\infty}^{i}=z_{0}^{i}+\int_{0}^{\infty} \dot{z}^{i}(t) \partial t \\
& =z^{i}(0)+\int_{\xi_{n}(0)}^{-\infty}\left(\partial z^{i} / \partial \xi_{n}\right) \partial \xi_{n}
\end{aligned}
$$

is well-defined since the integral converges.
Next, (1.5) and (1.6) imply

$$
\begin{aligned}
& \xi_{\mathrm{n}} \geq-\mathrm{C} / \rho, \quad \dot{\xi}_{\mathrm{n}} \leq-\mathrm{k}_{2} / \rho \\
\Rightarrow & \partial \xi_{\mathrm{n}} / \partial \mathrm{t} \leq \mathrm{k} \xi_{\mathrm{n}}, \quad \mathrm{k}>0 .
\end{aligned}
$$

Recalling that $\xi_{n}<0$, this differential inequality integrates to

$$
\xi_{n}(t) \leq-C e^{k t} .
$$

On the other hand, $\dot{\xi}_{\alpha}=O(1), \quad \alpha=1, \cdots, n-1, \quad$ and thus

$$
\left|\xi_{\alpha}(t)\right| \leq C t
$$

Finally,

$$
\begin{gathered}
d z^{i} / d z^{n}=\dot{z}^{i} / \dot{z}^{n}=\rho^{2} h^{i j} \xi_{j} / \rho^{2} h^{n j} \xi_{j} \\
=\frac{h^{i \alpha}\left(\xi_{\alpha} / \xi_{n}\right)+h^{i n}}{h^{n \alpha}\left(\xi_{\alpha} / \xi_{n}\right)+h^{n n}}
\end{gathered}
$$

By assumption $h^{\text {in }} \longrightarrow 0, h^{n n} \longrightarrow 1$, and by the estimates above, $\xi_{\alpha} / \xi_{n} \longrightarrow 0$. We conclude that $d z^{i} / d z^{n} \longrightarrow 0$.

## C. Curvature Asymptotics

The conformally compact metric $g$ of (1.2) is patently modelled on the Poincare metric on the unit ball in $\mathbb{R}^{n}$. The obvious desire is that much of the well understood hyperbolic geometry (and analysis) will be reflected in this more general setting. Indeed, this is the case with geodesic behaviour at infinity, as we have demonstrated in the last section. It is also true that the sectional curvatures of $g$ are negative when $\rho$ is small;
in fact there is a well defined limiting curvature tensor for which $\partial M$ is isotropic.
(1.10) Proposition: Let $\gamma(t)$ be a geodesic which approaches $p \in \partial M$. The sectional curvatures in any direction at $r(\mathrm{t})$ tend to $-\left(\partial \rho / \partial z^{n}\right)^{2}(\mathrm{p})$ as $\mathrm{t} \longrightarrow \infty$. Here, as usual, $\partial / \partial z^{n}$ is the $h$-unit normal to $\partial M$ at p.

Remark: Since $\rho_{n} \neq 0$ on $\partial M$, this proposition implies that the sectional curvatures are bounded between two negative constants in some collar neighbourhood of the boundary. Also, recall from the end of section $A$ that we may replace $\rho$ by $\phi \rho, h$ by $\phi^{2} h$, for any strictly positive function $\phi$ without altering $g$. The quantity $\partial \rho / \partial z^{n}$ is invariant under this transformation at $\rho=0$. Indeed, the new $\phi^{2} h$-unit normal is $\phi^{-1} \partial / \partial z^{n}$ and

$$
\phi^{-1} \partial_{z^{n}}(\phi \rho)=\rho_{\mathrm{n}}+\phi^{-1} \partial \phi / \partial \mathrm{z}^{\mathrm{n}} \rho=\rho_{\mathrm{n}} \quad \text { at } \quad \rho=0 .
$$

Proof: For the duration of the proof only, we use coordinates $\left(z^{1}, \cdots, z^{n}\right)$ which are normal with respect to $h$, for which $p$ corresponds to 0 , and such that $\partial / \partial z^{n}$ is orthogonal to $\partial M$ at $p$. Thus $\partial \rho / \partial z^{i}(p)=0, \quad i=$ $1, \cdots, \mathrm{n}-1, \quad \partial \rho / \partial \mathrm{z}^{\mathrm{n}}(\mathrm{p})>0$. and

$$
h_{i j}=\delta_{i j}+o(|z|)
$$

Set $u=\log \rho$, with subscripts on $u$ referring to $z^{i}$ derivatives, and denote by $\tilde{\Gamma}_{i j}^{k}$ the Christoffel symbols of $h_{i j}$. Notice that $\tilde{\Gamma}_{i j}^{k}=O(|z|)$. The Christoffel symbols of $\quad g_{i j}$ are

$$
\begin{equation*}
\Gamma_{i j}^{k}=\tilde{\Gamma}_{i j}^{k}-u_{i} \delta_{j}^{k}-u_{j} \delta_{i}^{k}+u_{e} h^{\ell k_{h}}{ }_{i j} \tag{1.11}
\end{equation*}
$$

Insert these into the formula expressing the components of the curvature tensor in terms of the $\Gamma^{\prime} s$ and their derivatives. A straightforward but messy calculation shows that

$$
\begin{gathered}
R_{i j i j}=\rho^{-2}\left\{\tilde{R}_{i j i j}+u_{i i}-2 u_{i j} \delta_{i j}+u_{j j}\right. \\
\left.+|\nabla u|^{2}\left(\delta_{i j}-1\right)+u_{i}^{2}-2 u_{i} u_{j} \delta_{i j}+u_{j}^{2}+o\left(|x|^{2}\right)\right\}
\end{gathered}
$$

Here $R$, $\tilde{R}$ are components of the curvature tensor for $g$, respectively $h$. Substitute

$$
u_{i}=\rho_{i} / \rho, \quad u_{i i}=\rho \rho_{i i} / \rho^{2}-\rho_{i}^{2} / \rho^{2}
$$

and assume $i \neq j$ to get

$$
R_{i j i j}=\rho^{-2} \tilde{R}_{i j i j}+\rho^{-4}\left\{\rho\left(\rho_{i i}+\rho_{j j}\right)-|\nabla \rho|^{2}+\rho^{2} O\left(|z|^{2}\right)\right\}
$$

Therefore, the sectional curvature in the $i-j$ plane is

$$
\frac{R_{i, j i j}}{\rho^{-4}\left(h_{i i} h_{j j}-h_{i j}^{2}\right)} \longrightarrow-|\nabla \rho|^{2}=-\left[\frac{\partial \rho}{\partial z^{n}}\right]^{2}
$$

as $\mathrm{t} \longrightarrow \infty$.
It is easy to follow through this proof using the metric $\rho^{-\sigma} h$, for any $\sigma>0$, instead. Indeed, the new $\Gamma_{i j}^{k}$ are obtained from those for $\sigma=2$ by replacing u with $\frac{1}{2} \sigma u$. The sectional curvature in the $i-j$ plane now becomes
$-\frac{\sigma^{2}}{4} \rho^{\sigma-2}|\nabla \rho|^{2}+\frac{\sigma}{2}\left(\frac{\sigma}{2}-1\right) \rho^{\sigma-2}\left(\rho_{i}^{2}+\rho_{j}^{2}\right)+\frac{\sigma}{2} \rho^{\sigma-1}\left(\rho_{i i}+\rho_{j j}\right)+O\left(\rho^{\sigma}\right)$.

This substantiates the claim of section $A$ that only when $\sigma=2$ is the curvature bounded between two negative constants near $\partial M$.

## D. Boundary Regularity

By their very definition, conformally compact manifolds have $C^{\infty}$ compactifications, but it is quite unclear to what extent this regularity is intrinsic to the metric g. In the present section we partially resolve this issue. Geometric regularity of an 'ideal' boundary is naturally gauged by the asymptotic properties of diverging geodesics. In this particular setting, let $p \in M$ and $\gamma(t)$ be a geodesic ray emanating from $p$ and tending toward the
boundary. Define

$$
\pi_{p}: w \longrightarrow \lim _{t \rightarrow \infty} \exp _{p}(t w)
$$

This maps a neighbourhood of $\boldsymbol{r}^{\prime}(0)$ in $\mathbf{S}_{\mathrm{p}} \mathrm{M}$ (the unit sphere in $T_{p} M$ ) to $\partial M$. Since $g$ has negative curvature bounded away from zero in the cone around $r(t)$, the results of Eberlein-O'Neill [15] imply that $\boldsymbol{\pi}_{\mathrm{p}}$ is a local homeomorphism. To impose an intrinsic smoothness structure on $\partial M$ it is sufficient to show that the composition maps $\pi_{q}^{-1} \pi_{p}: S_{p} M \longrightarrow S_{q} M$ all possess a certain fixed regularity. Anderson-Schoen [3] prove quite generally that each such map is $C^{\alpha}$, where the Hölder exponent $\alpha$ depends only on the upper and lower curvature bounds. Here, however, we are able to strengthen this considerably, but only with the strong proviso that the limiting sectional curvature function on $\partial M-\left(\partial \rho / \partial z^{n}\right)^{2}$ is constant in a neighbourhood of $\pi_{p}\left(r^{\prime}(0)\right)=r_{\infty}$. This bears marked resemblance to Fefferman's analogous result [16] for the Bergman metric of a strictly pseudoconvex domain, in which case the curvature is asymptotically constant along the whole boundary. I do not yet understand how essential this hypothesis is for conformally compact metrics.
(1.12) Proposition: Suppose, for the metric $g=\rho^{-2} h$, that $-\left(\partial \rho / \partial z^{n}\right)^{2}$ is constant on a neighbourhood of
$\pi_{p}\left(\gamma^{\prime}(0)\right)=\gamma_{\infty}$ in $\partial M$. Then $\pi_{p}$ maps a neighbourhood of $\gamma^{\prime}(0) \in S_{p} M$ diffeomorphically to a neighbourhood of $\gamma_{\infty}$ in $\partial M$.

Proof: The technique, inspired by [16], is to reparametrize the geodesic equations (1.3) so as to obtain a nonsingular system on a finite interval. It is convenient to choose coordinates so that $h^{n n}=1, h^{\alpha n}=0$. Here and in the following, we assume all Greek indices $\alpha, \beta, \gamma, \cdots$ vary between 1 and $n$, while $i, j, \ldots$ take the values $1, \cdots, \mathrm{n}$. By renormalizing we also assume $\rho(\gamma(0))=1$. Introduce new functions $v^{1}, \cdots, v^{n}$ by

$$
\begin{equation*}
\dot{z}^{\alpha}=\rho^{2} v^{\alpha}, \quad \dot{z}^{n}=\rho v^{n} \tag{1.13}
\end{equation*}
$$

so that

$$
\begin{equation*}
v^{\alpha}=h^{\alpha \beta} \xi_{\beta}, \quad v^{n}=\rho \xi_{n} . \tag{1.14}
\end{equation*}
$$

From the relationship (1.4) we have

$$
\rho^{2} h_{\alpha \beta} v^{\alpha} v^{\beta}+\left(v^{n}\right)^{2}=1
$$

By differentiating (1.14) and using this last equation, a bit of algebra shows that the $v^{i}$ satisfy

$$
\dot{\mathrm{v}}^{\alpha}=\rho^{2} h_{\beta \delta} \frac{\partial h^{\alpha \beta}}{\partial z^{\gamma}} \mathrm{v}^{\gamma} \mathrm{v}^{\delta}+\rho h_{\beta \gamma} \frac{\partial h^{\alpha \beta}}{\partial z^{n}} \mathrm{v}^{\gamma} \mathrm{v}^{\mathrm{n}}
$$

(1.15)

$$
\begin{aligned}
& -\frac{h^{\alpha \beta} \rho_{\beta}}{\rho}-\frac{\rho^{2}}{2} h_{\mu \sigma} h_{v \tau^{h}}{ }^{\alpha \beta} \frac{\partial h^{\mu v}}{\partial z^{\beta}} v^{\sigma} v^{\tau} \\
\dot{v}^{n}= & \rho \rho_{\alpha} v^{\alpha} v^{n}-\rho^{2} \rho_{n^{h}}{ }_{\alpha \beta} v^{\alpha} v^{\beta} \\
& -\frac{1}{2} \rho^{3} h_{\mu \alpha} h_{v \beta} \frac{\partial h^{\mu v}}{\partial z^{n}} v^{\alpha} v^{\beta} .
\end{aligned}
$$

Now introduce $z^{n}$ as the new parameter. The resulting equations of motion are

$$
\mathrm{dz} / \mathrm{z}^{\alpha} \mathrm{n}^{\mathrm{n}}=\rho \mathrm{v}^{\alpha} / \mathrm{v}^{\mathrm{n}}
$$

(1.16) $\quad \mathrm{dv}^{\alpha} / \mathrm{d} \mathrm{z}^{\mathrm{n}}=\rho h_{\beta \delta} \frac{\partial h^{\alpha \beta}}{\partial z^{\gamma}} \mathrm{v}^{\gamma} \mathrm{v}^{\delta} / \mathrm{v}^{\mathrm{n}}+\mathrm{h}_{\beta \gamma} \frac{\partial \mathrm{h}^{\alpha \beta}}{\partial \mathrm{z}^{\mathrm{n}}} \mathrm{v}^{\gamma}$

$$
-h^{\alpha \beta} \rho_{\beta} /\left(\rho^{2} v^{n}\right)-\frac{1}{2} \rho h_{\mu \sigma} h_{v \tau} h^{\alpha \beta} \frac{\partial h^{\mu v}}{\partial z^{\beta}} v^{\sigma} v^{\tau} / v^{n}
$$

$$
d v^{n} / d z^{n}=\rho_{\alpha} v^{\alpha}-\rho \rho_{n} h_{\alpha \beta^{2}} v^{\alpha} \beta / v^{n}
$$

$$
-\frac{1}{2} \rho^{2} h_{\mu \alpha} h_{v \beta} \frac{\partial h^{\mu v}}{\partial z^{n}} v^{\alpha} v^{\beta} / v^{n}
$$

These equations are nearly regular, the obstructions
being only the $v^{n}$ factors in the denominators and the third summand in the equation for $v^{\alpha}$. However, from (1.14) and the proof of (1.9) we have the a priori bound

$$
0>-c_{1} \geq v^{n} \geq-c_{2}
$$

where $c_{1}$, $c_{2}$ are independent of the particular geodesic. This resolves the former difficulty; as for the other one write $\rho(z)=a\left(z^{\prime}\right) z^{n}+b, \quad b=O\left(\left(z^{n}\right)^{2}\right)$. The coefficient function $a\left(z^{\prime}\right)$ is naturally associated with the metric g, being simply $\partial \rho /\left.\partial z^{n}\right|_{\rho=0}$. Hence $\rho_{\beta}=0\left(\rho^{2}\right)$ precisely when $a\left(z^{\prime}\right)$ is locally constant, and in this case $\rho_{\beta} / \rho^{2}$ is smooth.

Since we already know that solutions of (1.16) exist down to $z^{n}=0$, the standard theory for ordinary differential equations-in particular, smooth dependence on initial conditions-implies that the map

$$
\left(z^{\prime}(0), v(0)\right) \longrightarrow\left(z_{\infty}^{\prime}, v_{\infty}\right)
$$

from the initial values at $t=0$ to limiting values as $t \longrightarrow \infty$ is a $C^{\infty}$ diffeomorphism from the energy surface $|v(0)|=1$ to its image. Of course we need to show that

$$
\pi_{r(0)}: v(0) \longrightarrow z_{\infty}^{\prime}
$$

is a diffeomorphism. It has already been noted that $\pi_{\gamma(0)}$ is known to be a homeomorphism, so it suffices to prove that its differential is nonzero.

To this end a Jacobi field argument is natural. If $w \in T_{r(0)} M$ is orthogonal to $r^{\prime}(0)$, then essentially by definition

$$
d \pi_{r(0)}(w)=\lim _{t \rightarrow \infty} J_{w}(t)
$$

Here $J_{w}(t)$ is the Jacobi field along $r(t)$ such that $J_{W}(0)=0, J_{W}^{\prime}(0)=w$. It satisfies the equation

$$
\mathrm{J}_{\mathrm{w}}^{\prime \prime}+\mathrm{R}\left(\mathrm{~J}_{\mathrm{w}}, r^{\prime}\right) r^{\prime}=0
$$

By calculations very similar to those in section $C$, the components of $R$ and of the corresponding curvature operator $R_{0}$ for the metric of constant curvature $-a^{2}$ both blow up as $\rho^{-2}$. Furthermore, using strongly that $\rho_{\alpha}=$ $O\left(\rho^{2}\right)$, it is also the case that

$$
\left|R-R_{0}\right|=O(1)
$$

Hence, by the technique of asymptotic integration we conclude that $J_{w}$ grows at the same rate as the corresponding constant curvature Jacobi field $J_{w}^{0}$. With respect to the basis $\partial / \partial z^{1}, \cdots \partial / \partial z^{n}$, this latter field
tends to a finite nonzero limit as $t \longrightarrow \infty$. Thus $J_{w}$ does also, and the proof is complete.

## E. Laplace's Operator

This chapter concludes with a calculation to express the action of the Laplacian for the metric $g$ on differential forms. We use, now and for the rest of this paper, coordinates $\left(z^{1}, \cdots, z^{n}\right)$ with $z^{n}$ vanishing simply at $\partial M$, and $\partial / \partial z^{n}$ the $h$-unit normal there.

With the hindsight of experience, the metric on $\Lambda^{k_{M}}$

$$
\begin{equation*}
\left.\langle,\rangle_{g}=\rho^{2 k_{<}},\right\rangle_{h} \tag{1.17}
\end{equation*}
$$

is best regarded as a nondegenerate metric on singular k-forms. Introduce the space of sections

$$
\left\{\frac{\mathrm{d} \mathrm{z}^{1}}{\rho}, \cdots, \frac{\mathrm{dz}}{\rho}\right\}
$$

over $T^{*} M$. In the next chapter we shall define a new vector bundle for which these are nonsingular sections. At any rate, $k$-fold wedge products of these 1 -forms give a basis for $\Lambda^{k_{M}}$ over any point of the interior. Smooth combinations of these basis forms may be written

$$
\tilde{\omega}=\frac{1}{\rho^{k}} \omega
$$

where $\omega$ is a genuinely smooth form on $M . \quad$ If $\tilde{\eta}=\rho^{-k} \eta$, then

$$
\langle\tilde{\omega}, \tilde{\eta}\rangle_{\mathrm{g}}=\langle\omega, \eta\rangle_{\mathrm{h}} .
$$

Instead of studying the operation $\tilde{\omega} \longrightarrow \Delta_{g} \tilde{\omega}^{\sim}$ it is more convenient to study the induced operation on $\omega$ :

$$
\begin{equation*}
\Delta_{g}\left(\rho^{-k} \omega\right)=\rho^{-k} P \omega, \quad P=\rho^{k} \Delta_{g} \rho^{-k} \tag{1.18}
\end{equation*}
$$

Much of the 'hard' analysis later will focus on determining the mapping properties of a simpler operator which models $P$.

Now, since there is a factorization

$$
P=\left(\rho^{k} \mathrm{~d} \rho^{-k+1}\right)\left(\rho^{\mathrm{k}-1} \delta_{g} \rho^{-k}\right)+\left(\rho^{k} \delta_{g} \rho^{-k-1}\right)\left(\rho^{\mathrm{k}+1} \mathrm{~d} \rho^{-k}\right)
$$

it is easier to first calculate the conjugates of $d$ and $\delta_{g} . \quad$ First, then,

$$
\rho^{j+1} \mathrm{~d} \rho^{-\mathrm{j}} \omega=\rho^{\mathrm{j}+1}\left(\rho^{-\mathrm{j}} \mathrm{~d} \omega-\mathrm{j} \rho^{-\mathrm{j}-1} \mathrm{~d} \rho \Lambda \omega\right)
$$

(1.19)

$$
=\rho \mathrm{d} \omega-\mathrm{jd} \rho \wedge \omega
$$

As for the adjoint, use duality:

$$
\begin{aligned}
\int\left\langle\mathrm{d} \eta, \rho^{-j} \omega\right\rangle_{\mathrm{g}} \mathrm{dg} & =\int\langle\mathrm{d} \eta, \omega\rangle_{\mathrm{h}} \rho^{\mathrm{j}-\mathrm{n}} \mathrm{dh} \\
& =\int\left\langle\eta, \delta_{h}\left(\rho^{j-n_{\omega}}\right\rangle_{h} \mathrm{dh}\right. \\
& =\int\left\langle\eta, \delta_{h}\left(\rho^{j-n_{\omega}}\right)\right\rangle_{g} \rho^{2-2 j+n_{d g}} \\
& =\int\left\langle\eta, \rho^{-j+1}\left[\rho^{\mathrm{n}-\mathrm{j}+1} \delta_{h} \rho^{j-n_{\omega}} \omega\right\rangle_{g} \mathrm{dg}\right.
\end{aligned}
$$

for any j-form $\omega$ and j-1-form $\eta$. In other words

$$
\rho^{j-1} \delta_{g} \rho^{-j}=\rho^{n-j+1} \delta_{h} \rho^{j-n}
$$

Furthermore, for any metric

$$
\delta(\mathrm{f} \cdot \alpha)=\mathrm{f} \delta \alpha=\iota_{\nabla \mathrm{f}} \alpha
$$

Combining these last two formulae yields

$$
\begin{equation*}
\rho^{j-1} \delta_{g} \rho^{-j} \omega=\rho^{n-j+1}\left(\rho^{j-n} \delta_{h} \omega-(j-n) \rho^{j-n-1} \iota \quad \rho^{\omega}\right) \tag{1.20}
\end{equation*}
$$

$$
=\rho \delta_{h} \omega+(n-j) \iota_{\nabla} \rho^{\omega} .
$$

Here $\nabla \rho$ is the gradient of $\rho$ with respect to h. Finally, (1.19) and (1.20) together, along with a bit of computation, show that

$$
P \omega=\rho^{2} \Delta_{h} \omega+(2-k) \rho d \rho \wedge \delta_{h} \omega-k \rho \delta_{h}(d \rho \wedge \omega)
$$

$$
\begin{aligned}
& \text { (1.21) } \quad \begin{array}{l}
(n-k) \rho L_{\nabla \rho} \omega-2 \rho \iota \iota_{\nabla \rho} d \omega \\
\\
+(n-2 k) d \rho \wedge \iota_{\nabla \rho} \omega-k(n-k-1) \iota_{\nabla \rho}(d \rho) \omega \\
L_{\nabla \rho}=d \iota_{\nabla \rho}+\iota_{\nabla \rho} d \quad \text { is the Lie derivative. }
\end{array}
\end{aligned}
$$

Chapter 2. The Calculus of $V_{0}$ Pseudodifferential Operators

## A. Vo Vector Fields

The only distinguished submanifold of a manifold with boundary $M$ is its boundary. Related to this is the observation that the ring of differential operators Diff(M) has two natural subrings of geometric origin: Diff ${ }_{b}(M)$ and $\operatorname{Diff}_{\mathrm{O}}(\mathrm{M})$. The former is the space of operators which are sums of products of vector fields, unrestricted in the interior but required to lie tangent to the boundary-the so-called totally characteristic vector fields, the class of which is denoted $V_{b}$. The latter space is defined analogously using the $V_{0}$-vector fields, the ones vanishing at $\partial M$. Both $V_{b}$ and $V_{0}$ are Lie algebras under the usual bracket operation for vector fields.

Diff $f_{b}(M)$ and $\Psi_{b}(M)$, the related space of totally characteristic pseudodifferential operators, were introduced and studied by R. Melrose in [22]; Melrose-Mendoza [23] contains further developments. It is the purpose of this chapter to define and examine the ring $\Psi_{0}(M)$ of pseudodifferential operators generalizing Diffo(M). Much of the material here derives from information gleaned during conversations with Richard Melrose.

$$
\text { If } z=(y, x) \in \mathbb{R}^{n-1} \times \mathbb{R}^{+} \text {are coordinates of the }
$$

usual type, with $z^{n}(=x)$ vanishing on $\partial M$, then $V_{0}$ is generated as a $C^{\infty}(M)$-module by
(2.1)

$$
\left\{z^{n} \frac{\partial}{\partial z^{1}}, \cdots, z^{n} \frac{\partial}{\partial z^{n}}\right\} .
$$

Thus a typical element of $D_{i f f}^{m}(M)$ has an expression

$$
\begin{equation*}
P=\sum_{|\alpha| \leq m} a_{\alpha}(z)\left(z^{n} \frac{\partial}{\partial z}\right)^{\alpha} . \tag{2.2}
\end{equation*}
$$

The (modified) Laplacian of a conformally compact metric (1.21) lies in Diff ${ }_{0}^{2}(M)$.

Associated to $V_{0}$ is the group composed of those diffeomorphisms of $M$ which fix $\partial M$ pointwise. The expo-
 Any diffeomorphism induces a linear action on the tangent space of one of its fixed points. In particular, if $p \in$ $\partial M$ and $M_{p}$ is the inward pointing half of $T_{p} M$, a member of this class induces a linear transformation on $T_{p} M$ which preserves $M_{p}$. Such a map is of the form

$$
\begin{equation*}
(x, y) \longrightarrow(s x, y+x u) \tag{2.3}
\end{equation*}
$$

$$
s \in \mathbb{R}^{+}, u \in \mathbb{R}^{n-1} \simeq \partial M_{p}
$$

where now ( $x, y$ ) are linear coordinates on $M_{p}$ (with the order reversed from the previous paragraph). Let $G_{p}$ denote this linear group which is the semidirect product of $\mathbb{R}^{+}$with $\mathbb{R}^{n-1}$. Its composition and inverse laws are

$$
(s, u) \cdot(t, v)=(s t, t u+v)
$$

(2.4)

$$
(\mathrm{s}, \mathrm{u})^{-1}=\left(\mathrm{s}^{-1},-\mathrm{s}^{-1} \mathrm{u}\right)
$$

(2.3) should be recognized as the left action of $G_{p}$ on itself. In addition, the vector fields in (2.1) are the infinitesimal generators for this action. For future reference let us record the right action, along with its invariant vector fields

$$
(x, y) \longrightarrow(x, y) \cdot(s, u)^{-1}=\left(s^{-1} x, s^{-1}(y-u)\right)
$$

$$
\begin{equation*}
\left\{-x \frac{\partial}{\partial x}-y \cdot \frac{\partial}{\partial y}, \frac{\partial}{\partial y^{1}}, \cdots, \frac{\partial}{\partial y^{n-1}}\right\} \tag{2.5}
\end{equation*}
$$

It is interesting and significant to identify $G_{p}$ with the solvable subgroup $S=A N$ in the Iwasawa decomposition $G=A N K$ of the group of hyperbolic isometries:

$$
\mathrm{G}=\mathrm{SO}_{0}(1, \mathrm{n}), \quad \mathrm{A}=\mathbb{R}^{+}, \quad \mathrm{N}=\mathbb{R}^{\mathrm{n}-1}, \quad \mathrm{~K}=\mathrm{SO}(\mathrm{n})
$$

Furthermore, the $S$ invariant metric on $G / K \simeq M_{p}$ is hyperbolic. Thus, the introduction of the $V_{0}$ vector fields leads rather inexorably to constant negative curvature geometry.

In order to view $V_{O}$ somewhat more globally it is convenient to define the $V_{0}$ tangent bundle ${ }^{\circ} \mathrm{TM}$ as that bundle, the $C^{\infty}$ sections of which are precisely the $V_{0}$ vector fields

$$
\mathrm{v}_{\mathrm{O}}=\Gamma\left({ }^{\mathrm{o}} \mathrm{TM}\right)
$$

This way of constructing a bundle will be used several times later on, so we pause to elaborate briefly in this case. The individual fibre ${ }^{\circ} \mathrm{T}_{\mathrm{p}} \mathrm{M}$ is obtained by an equivalence relation on $V_{0}$ :

$$
\begin{aligned}
V, V^{\prime} \in V_{0} ; \quad V_{p} \sim V_{p}^{\prime} & \Leftrightarrow \\
V f=V^{\prime} f \quad \forall f \in C^{\infty}(M), & p \notin \partial M ;\left.d\left(V-V^{\prime}\right) f\right|_{p}=0 \\
& \forall f \in C^{\infty}(M), p \in \partial M .
\end{aligned}
$$

(2.1) exhibits a spanning set of sections with no rebalions; these determine a unique $C^{\infty}$ bundle structure on ${ }^{\circ}$ TM. The dual bundle ${ }^{\circ} \mathrm{T}^{*} \mathrm{M}$ is also quite useful. A basis of sections here, dual to the vector fields of (2.1), is

$$
\begin{equation*}
\left\{\frac{d z^{1}}{z^{n}}, \cdots, \frac{d z^{n}}{z^{n}}\right\}=\left\{\frac{d y}{x}, \frac{d x}{x}\right\} \tag{2.6}
\end{equation*}
$$

An operator in $\operatorname{Dif} f_{0}^{m}(M)$ has a well defined symbol, which is a homogeneous polynomial of degree $m$ on the fibres of ${ }^{\circ} \mathrm{T}^{*} \mathrm{M}$. If P is written as in (2.2), then

$$
\begin{equation*}
{ }^{o_{\sigma_{m}}(P)(z, \zeta)=} \sum_{|\alpha|=m} a_{\alpha}(z) \zeta^{\alpha} \tag{2.7}
\end{equation*}
$$

i.e. each $z^{n} \partial / \partial z^{i}$ is replaced by $\zeta^{i}$, and lower order terms are discarded. Invariance follows by continuity from the interior.

The operator $P$ is elliptic in this calculus if

$$
o_{\sigma_{m}}(\mathrm{P})(\mathrm{z}, \zeta)=0 \Leftrightarrow \zeta=0
$$

The Laplacian (1.21) is elliptic in this extended sense. Our goal in this paper is to prove Fredholm properties for this particular operator acting on certain weighted $L^{2}$ spaces. The general theory applies of course to any elliptic $V_{0}$ differential operator.

In the study of Diff $\mathrm{f}_{0}(\mathrm{M})$, as in the usual theory of standard differential operators, it is important to understand simpler models of operators in this class. In the interior of $M$ these are obtained by freezing the coefficients so as to obtain constant coefficient operators on each tangent space, also discarding all but the top order terms, exactly as usual. This last reduction is justifiable since $\partial_{z}^{\alpha}$ is 'stronger' when $|\alpha|=m$ than when
$|\alpha|<m$. At a point $p$ on the boundary, $P$ will be modelled more effectively by a certain $G_{p}$ invariant operator on $M_{p}$, which we call the normal operator $N_{p}(P)$. Its definition is suggested by the fact that near $z^{n}=0$

1, $z^{n} \frac{\partial}{\partial z^{i}},\left[z^{n} \frac{\partial}{\partial z^{i}}\right]\left[z^{n} \frac{\partial}{\partial z^{j}}\right], \quad$ and each $\left[z^{n} \frac{\partial}{\partial z}\right]^{\alpha}$
are all equally strong. However, any $\left(z^{n}\right)^{k}\left[\frac{\partial}{\partial z}\right]^{\alpha}$ with $k>|\alpha|$ is dominated by one of these operators, hence weaker.

Motivated by these observations we describe a procedure to define $N_{p}(P)$ when $P$ is differential. (2.43) further on contains a more general definition for pseudodifferential operators. We require first the notion of a normal fibration at $p \in \partial M$. This is a diffeomorphism from a neighbourhood $थ$ of $p$ to a neighbourhood $\|^{\prime}$ of $0 \in$ $T_{p} M$ such that

$$
f(p)=0, \quad f(थ) \subset M_{p}, \quad f(\partial थ) \subset \partial M_{p},\left.\quad f_{*}\right|_{0}=I d .
$$

Such a map is readily constructed using for example the geodesic flow of a metric for which $\partial M$ is totally geodesic. Then

$$
\begin{equation*}
N_{p}(P) u=\lim _{r \rightarrow 0}\left(R_{r}\right)^{*} f^{*} P\left(f^{-1}\right)^{*}\left(R_{1 / r}\right)^{*} u \tag{2.8}
\end{equation*}
$$

where $R_{r}$ is dilation by the factor $r$. Letting, for the moment, $\bar{z}$ denote the linear coordinates on $M_{p}$ induced by $z$ on $M$, then from (2.8) one may easily check

$$
\begin{equation*}
N_{p}\left(a(z)\left[z^{n} \frac{\partial}{\partial z}\right]^{\alpha}\right]=a(0)\left[\bar{z}^{n} \frac{\partial}{\partial \bar{z}}\right]^{\alpha}, \quad p=\{z=0\} \tag{2.9}
\end{equation*}
$$

(2.10)

$$
N_{p}(P \cdot Q)=N_{p}(P) \cdot N_{p}(Q), \quad P, Q \in \operatorname{Diff}_{0}(M)
$$

These formulae show that $N_{p}$ is a homomorphism from Diff ${ }_{O}(M)$ to the $G_{p}$ (left)-invariant operators on $M_{p}$.

## B. The Stretched Product Construction

Any attempt to construct a pseudodifferential inverse for an elliptic $P \in \operatorname{Diff}_{0}(M)$ must reconcile itself with the rather complete degeneracy of $P$ at $\partial M$. In particular, the Schwartz kernel of such an inverse must possess a singularity somewhat more severe than usual at the submanifold where the diagonal of $M \times M$ intersects the corner $\partial M \times \partial M$. In order the more thoroughly to display this singularity we shall define the appropriate class of kernels on a slightly larger manifold $M \times{ }_{0} M$, the $V_{0}$ stretched product of $M$ with itself. It is their natural abode; the additional singular behaviour is transferred to the geometry of $M \times{ }_{0} M$ and these kernels are as smooth as possible here.

We first examine the product $M \times M$ more closely. It has rather more natural structure than $M$ alone. Thus there are two natural hypersurfaces, the left and right boundaries

$$
\partial_{\ell}(M \times M)=\partial M \times M, \quad \partial_{r}(M \times M)=M \times \partial M .
$$

These are mapped surjectively to $M$ by the right and left projections $\pi_{r}$ and $\pi_{\ell}$, respectively, which send $M \times M$ onto the right and left factors. They intersect in the corner

$$
\partial_{\ell}(M \times M) \cap \partial_{r}(M \times M)=\partial M \times \partial M .
$$

Finally, the fixed point set of the involution $I$ which interchanges the two factors of $M$ is the diagonal $\Delta \iota$. This, in turn, is bounded by the diagonal of the corner


As noted earlier, the kernels of $V_{0}$-pseudodifferential operators should exhibit a singularity at $\partial \Delta c$ greater than their customary conormal one along the whole diagonal. This new behaviour is best exhibited by passing to the $V_{0}$ stretched product $M \times_{0} M$ where these kernels have a certain extension property. This manifold is defined by taking the real blow-up of $M \times M$ around $\partial \Delta c$, or to put it more familiarly, by introducing polar coordinates around this submanifold. Abstractly $M \times{ }_{0} M$ is formed from the disjoint union of $(\mathrm{M} \times \mathrm{M}) \backslash \partial \Delta \iota$ and $\mathrm{SN}_{++} \partial \Delta \iota$, the closed inward pointing sector of the spherical normal bundle to $\partial \Delta \iota$. [25] contains a proof that $M \times{ }_{0} M$ has a unique $C^{\infty}$ structure such that the blow-down map
(2.11)
$b: M \times{ }_{0} M \longrightarrow M \times M$
defined as the identity away from the new face, and as the bundle projection on $\mathrm{SN}_{++} \partial \Delta \iota$, is smooth and of rank $n-1$ (= dim $\partial \Delta \iota$ ) along this new face. We shall content ourselves by relying on coordinate descriptions.

Using the coordinates ( $x, y$ ) on the first factor of $M$ in the product, and an identical pair $(\tilde{x}, \tilde{y})$ on the second, then

$$
\Delta \iota=\{x=\tilde{x}, \quad y=\tilde{y}\}, \quad \partial \Delta \iota=\{x=\tilde{x}=0, \quad y=\tilde{y}\}
$$

Let $Y=y-\tilde{y}$, and define the singular coordinate system

$$
(\mathrm{R}, \omega, \tilde{\mathrm{y}})
$$

(2.12) $R=\left[x^{2}+|Y|^{2}+\tilde{x}^{2}\right]^{1 / 2}$,

$$
\begin{gathered}
\omega=R^{-1}(\mathrm{x}, \mathrm{Y}, \tilde{\mathrm{x}})=\left(\omega_{0}, \omega^{\prime}, \omega_{\mathrm{n}}\right) \\
R \geq 0, \omega \in \mathrm{~S}_{++}^{\mathrm{n}}=\left\{\omega \in \mathbb{R}^{\mathrm{n}+1}:|\omega|=1\right. \\
\\
\left.\omega_{0}, \omega_{\mathrm{n}} \geq 0\right\}
\end{gathered}
$$

so that

$$
x=R \omega_{0}, \quad \tilde{x}=R \omega_{n}
$$

(2.13)

$$
y=\tilde{y}+R \omega^{\prime}
$$

This system lifts to be $C^{\infty}$ on $M \times{ }_{0} M$. The new codimension one boundary, where $R=0$, is called the front face. The original two codimension one boundaries lift to the top and bottom faces:

$$
F=\{R=0\}, \quad T=\left\{\omega_{0}=0\right\}=b^{-1}(\{x=0\}) \backslash F
$$

(2.14)


It is often easier to picture $M \times{ }_{0} M$ instead as


Figure 3

The front face is the location where the analysis of the next chafer takes place. It has more structure than one might first suspect. First and foremost it is a quarter-sphere bundle over $\partial \Delta \iota$ with fibre $F_{p} \simeq S_{++}^{n}$ over
( $\mathrm{p}, \mathrm{p}$ ) , $\mathrm{p} \in \mathrm{O}_{\mathrm{M}}$. This follows either from (2.12) or its intrinsic description as $\mathrm{SN}_{++} \partial \Delta \iota$. Each $\mathrm{F}_{\mathrm{p}}$ also has a completely natural origin $O_{p}$, which is where $F_{p}$ intersects the lifted diagonal

$$
\begin{equation*}
\Delta \iota_{0}=\{\omega=(1 / \sqrt{2}, 0,1 / \sqrt{2}), \quad R, \tilde{y} \text { arbitrary }\} . \tag{2.15}
\end{equation*}
$$

The bundle of origins $0=F \cap \Delta \iota_{0}=\partial \Delta \iota_{O}$ is clearly searated from $\partial F$. In fact $\Delta \iota_{0}$ intersects $\partial\left(M \times{ }_{0} M\right)$ transversally, and only at the codimension one boundary. This remark indicates the fundamental gain incurred by passing to the stretched product, as we discuss in the next section.

The fibre $F_{p}$ admits two separate transitive $G_{p}$ actions in its interior, where $G_{p}$ is the group of (2.4). The groups

$$
G_{p}^{\ell}=G_{p} \times I d, \quad G_{p}^{r}=I d \times G_{p}
$$

act on the quarter space $M_{p} \times M_{p} \subset T_{(p, p)}(M \times M), \quad p \in \partial M$. Since $G_{p}$ fixes $\partial M_{p}$ they also fix $T(p, p) \partial \Delta \iota$, hence these actions descend first to the normal bundle

$$
\left.N_{++} \partial \Delta \iota\right|_{p}=T_{p}(M \times M)_{++} / T_{p} \partial \Delta \iota
$$

thence to its projectivization $\left.S N_{++} \partial \Delta \iota\right|_{p}=F_{p}$. Because of
the natural origin $O_{p}, F_{p}$ may be identified with $G_{p}$ either by way of the action of $G_{p}^{\ell}$ or that of $G_{p}^{r}$.
Through this identification these actions on $F_{p}$ corespond to the left or right action of $G_{p}$ on itself, and are intertwined by the (lift of the) involution I. Specifically, the map on $M_{p} \times M_{p}$
(2.16)

$$
\phi\left(\left(x_{1}, y_{1}\right),\left(x_{2}, y_{2}\right)\right)=\left(x_{1}, y_{1}\right) \cdot\left(x_{2}, y_{2}\right)^{-1}=\left[\frac{x_{1}}{x_{2}}, \frac{y_{1}-y_{2}}{x_{2}}\right]
$$

is invariant both under dilation and translations by elfments of $T(\partial \Delta \iota)$. Hence it diffeomorphically identifies the projectivization of (the interior of ) ( $\left.M_{p} \times M_{p}\right) / T_{p} \partial \Delta \iota=$ $F_{p}$ with $G_{p}$. Then $G_{p}^{\ell}$ acts as

$$
\begin{aligned}
& (t, v) \cdot\left(\left(x_{1}, y_{1}\right),\left(x_{2}, y_{2}\right)\right)=\left(\left(t x_{1}, y_{1}+x_{1} v\right),\left(x_{2}, y_{2}\right)\right) \\
& \xrightarrow{\phi}\left[t \frac{x_{1}}{x_{2}}, \frac{y_{1}-y_{2}}{x_{2}}+\frac{x_{1}}{x_{2}} v=(t, v) \cdot\left[\frac{x_{1}}{x_{2}}, \frac{y_{1}-y_{2}}{x_{2}}\right]\right.
\end{aligned}
$$

whereas for $G_{p}^{r}$

$$
\begin{aligned}
& (t, v) \cdot\left(\left(x_{1}, y_{1}\right),\left(x_{2}, y_{2}\right)\right)=\left(\left(x_{1}, y_{1}\right),\left(t x_{2}, y_{2}+x_{2} v\right)\right) \\
& \xrightarrow{\phi}\left[\frac{1}{t} \cdot \frac{x_{1}}{x_{2}}, \frac{1}{t} \frac{y_{1}-y_{2}}{x_{2}}-\frac{v}{t}\right]=\left[\frac{x_{1}}{x_{2}}, \frac{y_{1}-y_{2}}{x_{2}}\right] \cdot(t, v)^{-1}
\end{aligned}
$$

proving the claim.

Next we describe the infinitesimal $G_{p}^{\ell}$ action, which leads to an important perspective on the nature of $F_{p}$.
(2.17) Lemma: The vector fields $x \frac{\partial}{\partial x}, x \frac{\partial}{\partial y^{i}}$ on $M$ lift to the left factor of $M \times M$, thence to $M \times{ }_{0} M$. The restrictions of these lifts to $F$ lie tangent to each $F_{p}$; there they span the space of $G_{p}^{\ell}$-invariant vector fields. These restricted lifts are in $V_{0}$ with respect to (i.e. vanish on) $F_{p} \cap T$ and in $V_{b}$ with respect to (i.e. are tangent to) $F_{p} \cap B$. Finally, the quarter-sphere $F_{p}$ may be regarded as a ball $B^{n}$ blown up around a point on its boundary, with $F_{p} \cap B$ as the created face. The full space of 'mixed $V_{0}-V_{b}$ ' vector fields, generated by the $G_{p}^{\ell}$-invariant ones described above, is then the lift of $V_{0}\left(B^{n}\right)$ under this blow-up.


Figure 4

Remark: All assertions concerning tangency properties of various lifts in this statement are special cases of a
general property of the blow-up construction: if a manifold is blown-up around a submanifold, and $V$ is a vector field tangent to this submanifold, then it lifts on the blow-up to a vector field tangent to the new boundary, which is the pre-image of the submanifold.

Proof: Use the coordinates ( $x, y, \tilde{x}, \tilde{y}$ ) on $M \times M$ and $(R, \omega, \tilde{y})$, as in (2.12), on $M x_{0} M$. Then $x \partial_{x}, \quad x \partial_{y_{i}}$ lift to $M \times M$ to vector fields with the same expression. Next, if $b$ is the map of (2.11), then a brief computation shows that

$$
b_{*}\left(\omega_{0}^{2} R \partial_{R}+\omega_{0}\left\{\partial_{\omega_{0}}-\omega_{0} \sum_{j=0}^{n} \omega_{j} \partial_{\omega_{j}}\right\}\right)=x \partial_{x}
$$

(2.18)

$$
b_{*}\left(\omega_{0} \omega_{i} R \partial_{R}+\omega_{0}\left\{\partial_{\omega_{i}}-\omega_{i} \sum_{j=0}^{n} \omega_{j} \partial_{\omega_{j}}\right\}\right)=x \partial_{y_{i}}
$$

$$
i=1, \cdots, n-1
$$

Since $b$ is a diffeomorphism on the interior, and all these vector fields are smooth, the vector fields in parentheses here are the unique smooth lifts of $\quad x \partial_{x}, \quad x \partial_{y_{i}}$.

Their restrictions to $F$ are

$$
x=\omega_{0}\left\{\partial_{\omega_{0}}-\omega_{0} \sum_{j=0}^{n} \omega_{j} \partial_{\omega_{j}}\right\}
$$

(2.19)

$$
Y_{i}=\omega_{0}\left\{\partial_{\omega_{i}}-\omega_{i} \sum_{j=0}^{n} \omega_{j} \partial_{\omega_{j}}\right\}, \quad i=1, \cdots, n-1
$$

These are patently orthogonal to $\sum \omega_{j} \partial_{\omega_{j}}$, hence tangent to F. Furthermore, since they contain no $\partial_{y}^{\sim}$ terms, they are also tangent to each $F_{p}$. The factor $\omega_{0}$ ensures their vanishing at $F_{p} \cap T$ where $\omega_{0}=0$. On the other hand,

$$
X \cdot \partial_{\omega_{n}}=Y_{i} \cdot \partial_{\omega_{n}}=0
$$

at $F_{p} \cap B$, where $\omega_{n}=0$ and $\partial_{\omega_{n}}$ is the normal, so they are tangent here, as claimed.

It is by no means obvious that $X$ and $Y_{i}$ of (2.19) are $G_{p}^{\ell}$ invariant. It is convenient therefore to use the projective coordinates of (2.16)

$$
\begin{equation*}
s=x / \tilde{x}, \quad u=\frac{\tilde{y}-\tilde{y}}{\tilde{x}} \quad(x, y),(\tilde{x}, \tilde{y}) \in M_{p} \tag{2.20}
\end{equation*}
$$

These are nonsingular on the interior of $F_{p}$, and even on neighbourhoods of $F_{p} \cap T$ bounded away from $F_{p} \cap T \cap B$. Notice that

$$
\mathrm{s}=\frac{\omega_{0}}{\omega_{\mathrm{n}}}, \quad \mathrm{u}=\frac{\omega^{\prime}}{\omega_{\mathrm{n}}} .
$$

Another short calculation now shows that

$$
\begin{equation*}
X=s \partial_{s}, \quad Y_{i}=s \partial_{u_{i}} \tag{2.21}
\end{equation*}
$$

from which $G_{p}^{\ell}$ invariance is obvious.
As for the last assertions of (2.17), first observe that the blow-up of a ball around a point $q$ in its boundary is smoothly equivalent to a quarter-sphere. Denote the blow-down map by

$$
\beta: F_{p} \simeq S_{++}^{n} \longrightarrow B^{n}
$$

It carries $F_{p} \backslash F_{p} \cap B$ diffeomorphically to $B^{n} \backslash\{q\}$, and sends $F_{p} \cap B$ to $q$. We need to verify that the lifts of $V_{0}\left(B^{n}\right)$ vector fields vanish at $F_{p} \cap T$ and lie tangent to $F_{p} \cap B$. The vanishing is obvious from the bijectivity of $\beta$ near $F_{p} \cap T$; the tangency is proved by essentially the same computation as goes into (2.18), which we leave to the reader to check. This completes the proof.
(2.22) Definition: $E$ is the vector bundle over $F_{p}$ whose full space of $C^{\infty}$ sections are the vector fields vanishing at $F_{p} \cap T$ and tangent to $F_{p} \cap B$ described above. We
have shown that

$$
\beta^{*}\left(\Gamma\left({ }^{\mathrm{o}} \mathrm{~T}\left(\mathrm{~B}^{\mathrm{n}}\right)\right)\right)=\Gamma(\mathrm{E}) .
$$

Our description of the basic natural baggage carried by $M \times{ }_{0} M$ concludes with a discussion of certain vector bundles we shall be using. These are the $V_{0}$ density bundles and the $V_{0} k$-form bundles.

Let us commence with the density bundles, as they are somewhat simpler to understand. Recall the Riemannian density dg of (1.2), which grows as $\rho^{-n}$ near $\partial M$. Define then $\Gamma_{0}(M)$ to be the space of all $C^{\infty}$ densities $v$ on the interior of $M$ such that $\rho^{n} \cdot v$ extends smoothly to the closure. $\Gamma_{0}(M)$ is the space of sections of a $C^{\infty}$ line bundle over $M$, and a typical element may be written in coordinates as

$$
v=h(x, y)\left|x^{-n} d x d y\right| \quad h \in C^{\infty}(M)
$$

This generalizes naturally to a manifold $X$ with (codimension one) boundary components $X_{1}, \cdots, X_{N}$, which have defining functions $\rho_{1}, \cdots, \rho_{N}$, respectively. Let

$$
\begin{align*}
& \Gamma_{0}(X)=\{\text { smooth densities } v \text { on } \dot{X} \text { such }  \tag{2.23}\\
& \left.\quad \text { that } \rho_{1}^{n} \cdots \rho_{N}^{n} v \text { extends smoothly to } X\right\} .
\end{align*}
$$

The connection between these densities and $V_{0}$-geometry is illustrated by
(2.24) Lemma: Densities in $\Gamma_{0}(M \times M)$ lift to densities in $\Gamma_{0}\left(M \times{ }_{0}{ }^{M}\right)$ under the usual blow-down map. In fact, the line bundle for which the latter space is the space of sections is the pullback under $b$ of the line bundle corresponding to the former space.

Proof: Away from $F$ and $\partial \Delta \iota$ there is nothing to prove since $b$ is a diffeomorphism there. At $F$ use the coordinates ( $\mathrm{R}, \omega, \tilde{y}$ ) of (2.12):

$$
\left|\frac{d x d y d \tilde{x} d \tilde{y}}{x^{n} x^{n}}\right|=\left|\frac{R^{n} d R d \omega d \tilde{y}}{\left(R \omega_{0}\right)^{n}\left(R \omega_{n}\right)^{n}}\right|=\left|\frac{d R d \omega d \tilde{y}}{R^{n} \omega_{0}^{n} \omega_{n}^{n}}\right| .
$$

The left term spans $\Gamma_{0}(M \times M)$ since $x$ and $\tilde{x}$ are defining functions for the left and right boundary components, whereas the term on the right generates $\Gamma_{0}\left(M_{0} M\right)$ as $F$, T, B are defined by $R, \omega_{0}, \omega_{n}$, respectively. This proves the lemma.

One may also consider the powers of either of these density bundles; they too correspond under pullback. The half-densities are the ones we use later. Both $\Gamma_{0}^{1 / 2}(\mathrm{M} \times \mathrm{M})$ and $\Gamma_{0}^{1 / 2}\left(\mathrm{M} \times_{0} \mathrm{M}\right)$ will commonly be denoted simply $\Gamma_{0}^{1 / 2}$; the base space should be clear from the context.

Let us turn now to the k-form bundles; consider first the $V_{0}$ cotangent bundle ${ }^{\circ} \mathrm{T}^{*} \mathrm{M}$. By pulling it back to the left factor of $M \times M$, and then to $M \times{ }_{0} M$, one obtains a bundle ${ }^{o} \Lambda_{\ell}^{1}$ over the stretched product. The subscript $\ell$ signifies the involvement of the left factor of $M \times M$. By functoriality the pullback of $\Lambda^{k}\left({ }^{0} T^{*} M\right)$ as above coincides with $\Lambda^{k}\left({ }^{0} \Lambda_{\ell}^{1}\right)$. Call this new bundle ${ }^{\mathbf{o}} \Lambda_{\ell}^{k}$. Using the right factor of the product, ${ }^{0} \Lambda_{r}^{k}$ is also defined. Their spaces of sections are denoted ${ }^{0} \Omega_{\ell}^{k},{ }^{0} \Omega_{r}$.

The blow-down $b$ carries $F_{p}$ to the single point ( $p, p$ ) $\in \partial \Delta \iota$. Hence ${ }^{o} \Lambda_{\ell}^{k}$, for example, restricts to a canonically trivial bundle over $F_{p}$. To identify this restriction with a more familiar object, compute the lifts of $\frac{d x}{x}, \frac{d y_{i}}{x}$ on $M$ to $M x_{0} M$. Using the coordinate system (s,u,x,y), with (s,u) as in (2.20), which is nonsingular away from the bottom face $B$,

$$
\frac{d x}{x}=\frac{d s}{s}+\frac{d \tilde{x}}{\tilde{x}}, \quad \frac{d y_{i}}{x}=\frac{d \tilde{y}_{i}}{s \cdot \tilde{x}}+\frac{u_{i}}{s} \frac{d \tilde{x}}{\tilde{x}}+\frac{d u_{i}}{s} .
$$

We define their 'formal pullbacks' to $F_{p}$

$$
\begin{equation*}
\frac{d x}{x} \longrightarrow \frac{d s}{s}, \frac{d y_{i}}{x} \longrightarrow \frac{d u_{i}}{s} \tag{2.25}
\end{equation*}
$$

(2.26) Lemma: The correspondence (2.25) induces an isomorphism

$$
\left.{ }^{o} \Lambda_{\ell}^{\mathrm{k}}\right|_{\mathrm{F}_{\mathrm{p}}} \simeq \Lambda^{\mathrm{k}}\left(\mathrm{E}^{*}\right)
$$

where $E^{*}$ is the dual of the bundle $E$ described in (2.22).

Proof: $s^{-1} d s, s^{-1} d_{i}$ are dual to $s \partial_{s}, s \partial_{u_{i}}$, which by (2.17) generate the bundle E. In addition, we showed there that $E$ is a smooth bundle over the closed face $F_{p}$, hence $E^{*}$ also extends smoothly to the closure. The rest of the proof is obvious.

Remark: We shall prove a slight extension of this Lemma in (3.3). Also, the rules (2.25) apply to the $V_{0}$ density bundles since densities are simply 'absolute values' of n-forms. So analogously

$$
\left.\Gamma_{0}\left(\mathrm{M} \times{ }_{0} \mathrm{M}\right)\right|_{\mathrm{F}_{\mathrm{p}}} \simeq \Gamma_{0}\left(\mathrm{~F}_{\mathrm{p}}\right)
$$

## C. Vo Kernels

The temperate distributions on a manifold with boundary $M$ are those with specific extension properties across this boundary. If $\tilde{M}$ is an open region containing $M$ in its interior, then distributions on $\tilde{M}$ with support in $M$ constitute the space of distributions dual to $C^{\infty}(M)$. On the other hand, restrictions of arbitrary distributions on
$\tilde{M}$ to the interior of $M$ lie in the space dual to $\dot{C}^{\infty}(M)$-functions vanishing to infinite order on $\partial M$. The resulting spaces on $M$, of supported and extendible distributions, respectively, are quite similar but have slightly different extension properties nonetheless.

On a manifold-with-corner, for example $M \times M$, there are various other possibilities. For instance, an ordinary pseudodifferential operator on $M$ has a Schwartz kernel $k$ on $M \times M$ which has a conormal singularity along the diagoral and is $C^{\infty}$ elsewhere. In fact $k$ then extends to be conormal along the 'longer' diagonal of $\tilde{M} \times \tilde{M}$, and this in some sense fixes its singularity at $\partial \Delta \iota$. There are, however, other quite reasonable extension properties which $k$ could be required to have. To formulate one of these, recall that a manifold with boundary can be doubled across its boundary. In particular, we may form the double [ $M \times{ }_{0} M^{2}$ of the stretched product across its front face.


Figure 5

The lifted diagonal joins with its double, resulting in a boundary-less submanifold contained in the interior of $\left[\mathrm{M} \times \mathrm{O}_{\mathrm{M}}\right]^{2}$.
(2.27) Definition: $K_{0}^{m}\left(M ; \Gamma_{0}^{1 / 2}\right)$ is the space of distributional sections of $\Gamma_{0}^{1 / 2}\left(M \times{ }_{0} M\right)$ which extend to $\left[M \times{ }_{0} M^{2}\right.$ as distributions conormal of order $m$ (see [19]) along the extended diagonal (tensored with $C^{\infty}$ sections of a smooth extension of $\Gamma_{0}^{1 / 2}$ ) and vanishing to all orders on $\partial\left(\left[M \times{ }_{0}\right]^{2}\right)$.
$\Psi_{0}^{m}\left(M ; \Gamma_{0}^{1 / 2}\right)$, the space of $V_{0}$ pseudodifferential operators, is the collection of operators with Schwartz kernels on $M \times M$ lifting to elements of $K_{0}^{m}\left(M ; \Gamma_{0}^{1 / 2}\right)$.

Remark: If $A \in \Psi_{0}^{m}$, then we shall say that its kernel $\kappa(A)$ belongs to $K_{0}^{m}$. Notice also that the pushforward of $\kappa \in K_{0}^{m}$ is well defined: the blow-down map $b$ is a diffeomorphism away from $F$, and the intersection of $\partial \Delta \iota$ with $F$ is transversal so by wave front set considerations this operation makes sense.

Let us examine how $A \in \Psi_{0}^{m}\left(M ; \Gamma_{0}^{1 / 2}\right)$ acts. Set

$$
\mu=\left|\frac{\mathrm{dx} \mathrm{dy}}{\mathrm{x}^{\mathrm{n}}}\right|^{1 / 2}
$$

(2.28)

$$
r=\left|\frac{d s d u d \tilde{x} d \tilde{y}}{s^{n} x^{n}}\right|^{1 / 2}=\left|\frac{d x d y d s d u}{s \cdot x^{n}}\right|^{1 / 2}
$$

where $s$ and $u$ are defined in (2.20). The full system ( $s, u, \tilde{x}, \tilde{y}$ ) are regular except along $B$, and shall frequently be used. Then

$$
\begin{equation*}
\kappa(A)=k(s, u, \tilde{x}, \tilde{y}) \cdot \gamma \tag{2.29}
\end{equation*}
$$

where $k$ is conormal of order $m$ on $\{s=1, u=0\}$, and decreases rapidly as $s \longrightarrow 0, s \longrightarrow \infty$, or $|u| \longrightarrow \infty$. On the half density $f \cdot \mu$, we have

$$
\begin{equation*}
(A f)(x, y)=\left\{\int k\left(s, u, \frac{x}{s}, y-\frac{x}{s} u\right) f\left(\frac{x}{s}, y-\frac{x}{s} u\right) \frac{d s}{s} d u\right\} \cdot \mu \tag{2.30}
\end{equation*}
$$

since $\tilde{x}=s^{-1} x, \tilde{y}=y-s^{-1} x u$. In particular, the identity is represented by

$$
\begin{equation*}
\kappa(\mathrm{Id})=\delta(\mathrm{s}-1) \delta(\mathrm{u}) \cdot \gamma \in \mathrm{K}_{0}^{0}\left(\mathrm{M} ; \Gamma_{0}^{1 / 2}\right) \tag{2.31}
\end{equation*}
$$

We record another representation of this action, now using the coordinates $(x, y, t, v), t=\tilde{x} / x, v=(\tilde{y}-y) / x$. Here the half density $r$ equals

$$
\left|\frac{d x d y d t d v}{x^{n} t^{n}}\right|^{1 / 2}
$$

and

$$
\begin{equation*}
(A f)(x, y)=\left\{\int k^{\prime}(x, y, t, v) f(x t, y+x v) \frac{d t}{t^{n}} d v\right\} \cdot \mu \tag{2.32}
\end{equation*}
$$

if $k(A)=k(x, y, t, v)$ in these coordinates. Although (2.32) appears both more natural and neater than (2.30), the first expression is better for our purposes. The reason is that we are frequently interested in the behaviour of $A f$ as $x \longrightarrow 0$, but (t,v) are singular along the top face $T$, hence less convenient.

It is straightforward to extend these definitions to include operators acting on the k-form bundles of the last section. Thus

$$
\begin{align*}
A & \in \Psi_{0}^{\mathrm{m}}\left(\mathrm{M} ;{ }^{\mathrm{o}} \Lambda^{\mathrm{k}} \otimes \Gamma_{0}^{1 / 2}\right)  \tag{2.33}\\
& \Leftrightarrow \kappa(\mathrm{A}) \in \mathrm{K}_{0}^{\mathrm{m}}\left(\mathrm{M} ; \Gamma_{0}^{1 / 2}\right) \otimes \operatorname{Hom}\left({ }^{\mathrm{o}} \Lambda_{\mathrm{r}}^{\mathrm{k}},{ }^{\mathrm{o}} \Lambda_{\ell}^{\mathrm{k}}\right) .
\end{align*}
$$

Unfortunately the class of operators we have defined is not large enough to contain inverses of elliptic $V_{0}$ differential operators. For if $E$ is to invert $P \in$ Diff $\mathrm{O}_{\mathrm{O}}^{\mathrm{m}}$, it must somehow incorporate information about the asymptotic behaviour of (formal) solutions to $P u=0$ near the boundary. The kernels we have already introduced vanish to infinite order at the relevant boundary components $T$ and $B$.

We remedy this shortcoming as follows. Let $T$ and $B$ denote also the extended top and bottom faces in $\left[M \times{ }_{0}\right]^{2}$, with $\rho_{\mathrm{T}}$ and $\rho_{\mathrm{B}}$ their corresponding defining functions. Suppose $V_{b}$ is the space of vector fields on this double which are tangent to both boundary components; set

$$
\begin{aligned}
& A^{a, b}= \\
& \left\{u \in \mathscr{D} \cdot\left(\left[M \times 0^{M}\right]^{2}\right): V_{1} \cdots v_{j} u \in \rho_{T}^{a} \rho_{B}^{b}\left(\log \rho_{T} \log \rho_{B}\right)^{N_{L} L^{\infty}}\right. \\
& v_{i} \in V_{b}, \quad i=1, \cdots, j, \text { for all } \quad j \quad \text { and some } N \\
& \text { independent of } j\}
\end{aligned}
$$

This has an invariantly defined subspace of polyhomogeneous elements

$$
A_{\mathrm{phg}}^{\mathrm{a}, \mathrm{~b}}=\left\{u \in A^{\mathrm{a}, \mathrm{~b}}: u \sim \sum_{\mathrm{i}=0}^{\infty} \sum_{j=0}^{N_{i}} u_{i j} \rho_{\mathrm{T}}^{\mathrm{a+i}}\left(\log \rho_{\mathrm{T}}\right)^{j}\right.
$$

near $T$,
(2.34)

$$
\left.u \sim \sum_{i=0}^{\infty} \sum_{j=0}^{N N_{i}^{\prime}} u_{i j}^{\prime} \rho_{T}^{b+i}\left(\log \rho_{B}\right)^{j} \quad \text { near } \quad B\right\}
$$

Here $u_{i j}$ and $u_{i j}$ are $C^{\infty}$ and depend only on the 'tangential' variables at the appropriate boundary.
(2.35) Definition:

ii) $K_{0}^{\mathrm{m}, \mathrm{a}, \mathrm{b}}\left(\mathrm{M} ; \Gamma_{0}^{1 / 2}\right)=\mathrm{K}_{0}^{\mathrm{m}}\left(\mathrm{M} ; \Gamma_{0}^{1 / 2}\right)+\mathrm{K}_{0}^{-\infty, \mathrm{a}, \mathrm{b}}\left(\mathrm{M} ; \Gamma_{0}^{1 / 2}\right)$.
iii) $A \in \Psi_{0}^{m, a, b} \Leftrightarrow \kappa(A) \in K_{0}^{m, a, b}$.

The kernels in $i$ ) are just sections of $\Gamma_{0}^{1 / 2}\left(M \times{ }_{0}^{M}\right)$ which are $C^{\infty}$ in the interior and up to the front face, and conormal with classical expansions at $T$ and $B$.

The extension of these last definitions to operators acting on differential forms is slightly arcane since we need to consider forms with tangential and normal components vanishing with different rates at the boundary. Some additional structure is needed so that these components are well-defined; for the present problem we have the conformal class of the metric $h$ with which to work. Using this, the ordinary form bundle splits at $\partial M$ :

$$
\left.\Lambda^{k}(M)\right|_{\partial M}=\Lambda_{t}^{k} \oplus \Lambda_{n}^{k}
$$

Similarly, from the formal pullback (2.25) we have

$$
\left.{ }^{\mathrm{o}} \Lambda^{\mathrm{k}}(\mathrm{M})\right|_{\partial \mathrm{M}}={ }^{\mathrm{o}} \Lambda_{\mathrm{t}}^{\mathrm{k}} \oplus{ }^{\mathrm{o}} \Lambda_{\mathrm{n}}^{\mathrm{k}}
$$

Next, by virtue of the metric $h$ again, we may actually assume that such a splitting is defined in a neighbourhood of the boundary. This is less natural, but the important point is that in any two such splittings the summands agree to first order at the boundary. Another way to phrase this is to say that the conformal class distinguishes the collection of coordinate systems ( $y, x$ ) for which $x=0$ on the boundary and $\partial / \partial x$ is orthogonal to $\partial M$. Henceforth only such coordinates will be used, and we assume that a splitting of the $V_{0}$ form bundle is given.

The real purpose of these machinations is so that we may attach an invariant sense to the requirement that the components of $\omega=\omega_{t}+\omega_{n}$ vanish at different rates at $\partial \mathrm{M}$ :

$$
\omega_{\mathrm{t}}=0\left(\mathrm{x}^{\mathrm{a}}\right) \quad \omega_{\mathrm{n}}=0\left(\mathrm{x}^{\mathrm{b}}\right)
$$

If $\omega_{t}, \omega_{n}$ are computed with respect to the coordinates ( $y, x$ ) and are seen to vanish at these rates, and then recomputed with respect to new coordinates ( $\bar{y}, \bar{x}$ ), these rates might be intermingled. If, however, both coordinate systems are of the special type above, then from the fact that

$$
\frac{\partial y^{i}}{\partial \bar{x}}=O(\bar{x}), \quad \frac{\partial \mathrm{x}}{\partial \overline{\mathrm{y}}^{i}}=O(\overline{\mathrm{x}})
$$

it is easy to see that the new components $\bar{\omega}_{t}$, $\bar{\omega}_{t}$ vanish at the rates min(a,b+1), min(a+1,b), respectively. Hence, these rates are well defined provided we assume $b-1 \leq a \leq b+1$. Similarly, it is likewise invariant to require that

$$
\omega_{\mathrm{t}} \in A_{\mathrm{pgh}}^{\mathrm{a}}, \quad \omega_{\mathrm{n}} \in A_{\mathrm{phg}}^{\mathrm{b}}
$$

if $|a-b| \leq 1$. The spaces here are simply those of (2.34) when only one boundary is present, and of course are extended to contain vector-valued elements.

Now, using all of these technical hypotheses, we assume that both ${ }^{0} \Lambda_{\ell}^{k}$ and ${ }^{0} \Lambda_{r}^{k}$ split in a neighbourhood of $F$ in $M \times{ }_{0} M$ into their tangential and normal subbundles. Thus any $G \in \operatorname{Hom}\left({ }^{\mathbf{o}} \Lambda_{\ell}^{\mathbf{k}},{ }_{0} \Lambda_{r}^{k}\right)$ may be associated to a matrix

$$
G=\left[\begin{array}{ll}
G_{t t} & G_{n t}  \tag{2.36}\\
G_{t n} & G_{n n}
\end{array}\right]
$$

where

$$
G_{t t}:\left({ }^{\mathrm{o}} \Lambda_{\mathrm{r}}^{\mathrm{k}}\right)_{\mathrm{t}} \longrightarrow\left({ }^{\mathrm{o}} \Lambda_{\ell}^{\mathrm{k}}\right)_{\mathrm{t}} \quad \mathrm{G}_{\mathrm{nt}}:\left({ }^{\mathrm{o}} \Lambda_{\mathrm{r}}^{\mathrm{k}}\right)_{\mathrm{n}} \longrightarrow\left({ }^{\mathrm{o}} \Lambda_{\ell}^{\mathrm{k}}\right)_{\mathrm{t}}
$$

$$
G_{\mathrm{tn}}:\left({ }^{\mathrm{o}} \Lambda_{\mathrm{r}}^{\mathrm{k}}\right)_{\mathrm{t}} \longrightarrow\left({ }^{\mathrm{o}} \Lambda_{\ell}^{\mathrm{k}}\right)_{\mathrm{n}} \quad \mathrm{G}_{\mathrm{nn}}:\left({ }^{\mathrm{o}} \Lambda_{\mathrm{r}}^{\mathrm{k}}\right)_{\mathrm{n}} \longrightarrow\left({ }^{\mathrm{o}} \Lambda_{\ell}^{\mathrm{k}}\right)_{\mathrm{n}}
$$

Thus, for $\omega=\omega_{t}+\omega_{n}, \quad G \omega=(G \omega)_{t}+(G \omega)_{n}$ where

$$
\begin{aligned}
& (G \omega)_{t}=G t_{t} \omega_{t}+G_{n t} \omega_{n} \\
& (G \omega)_{n}=G_{t n} \omega_{t}+G_{n n} \omega_{n} .
\end{aligned}
$$

The $G_{i j}, i, j=t, n$, are also called the components of G .

These byzantine preliminaries allow us finally to define kernels generalizing those of (2.35). The idea is simply to let each $G_{i j}$, when expressed as a matrix, contain entries in $\alpha_{p h g}^{a, b}$. More precisely, suppose $\sigma, T$ are two by two real matrices

$$
\sigma=\left[\begin{array}{ll}
\sigma_{\mathrm{tt}} & \sigma_{\mathrm{nt}} \\
\sigma_{\mathrm{tn}} & \sigma_{\mathrm{nn}}
\end{array}\right] \quad \tau=\left[\begin{array}{cc}
\tau_{\mathrm{tt}} & \tau_{\mathrm{nt}} \\
\tau_{\mathrm{tn}} & \tau_{\mathrm{nn}}
\end{array}\right]
$$

(2.37) Definition:
i) $K_{0}^{-\infty, \sigma, \tau}\left(M ;{ }^{0}{ }^{\mathrm{k}} \otimes \Gamma_{0}^{1 / 2}\right)=\left\{G=\left(G_{i j}\right), i, j=t, n\right.$,

$$
G_{i j} \in K_{0}^{-\infty, \sigma_{i j}, \tau_{i j}}\left(M, \Gamma_{0}^{1,2}\right) \otimes \operatorname{Hom}\left(\left({ }^{o} \Lambda_{r}^{k}\right)_{i},\left({ }^{o} \Lambda_{\ell}^{k}\right)_{j}\right\}
$$

ii) $K_{0}^{m, \sigma, T}=K_{0}^{m}+K_{0}^{-\infty, \sigma, \tau}$.
iii) $A \in \Psi_{0}^{m, \sigma, T} \Leftrightarrow \kappa(A) \in K_{0}^{m, \sigma, \tau}$.

As before, these spaces are well defined only when

$$
\begin{aligned}
& \left|\min \left(\sigma_{t t}, \sigma_{n t}\right)-\min \left(\sigma_{t n}, \sigma_{n n}\right)\right| \leq 1 \\
& \left|\min \left(\tau_{t t}, \tau_{n t}\right)-\min \left(\tau_{t n}, \tau_{n n}\right)\right| \leq 1 .
\end{aligned}
$$

To conclude this treatment of $V_{0}$ pseudodifferential operators, we discuss their composition properties (esecoaly with $V_{0}$ differential operators) and the symbol isomorphism. The following results will only be stated and proved for scalar operators, but obviously still hold for operators on forms, by, for example, referring to a basis.
(2.38) Proposition: $\Psi_{0}^{*, a, b}$ is a two-sided Diff ${ }_{0}^{*}$ module

$$
\begin{aligned}
& \operatorname{Diff}_{0}^{\mu} \cdot \Psi_{0}^{\mathrm{m}, \mathrm{a}, \mathrm{~b}} \subset \Psi_{0}^{\mathrm{m}+\mu, \mathrm{a}, \mathrm{~b}} \\
& \Psi_{0}^{\mathrm{m}, \mathrm{a}, \mathrm{~b}} \cdot \operatorname{Diff}_{0}^{\mu} \subset \Psi_{0}^{\mathrm{m}+\mu, \mathrm{a}, \mathrm{~b}}
\end{aligned}
$$

Proof: It suffices first of all only to prove the first inclusion, for the other follows by taking adjoint (note that adjoint of elements in $\Psi_{0}^{m, a, b}$ lie in $\Psi_{0}^{m, b, a}$ ). Furthermore it also suffices to take $\mu=1$ and then iterate. Consider then the vector fields $x \partial_{x}, x_{y_{i}} . B y$
(2.17) these lift to $M x_{0} M$ to vector fields tangent to all boundaries, whence they preserve the conormal spaces $A^{\mathrm{a}, \mathrm{b}}$. Furthermore, using the coordinates (s,u) of (2.20) we showed there that

$$
b_{*}\left(s \partial_{s}\right)=x \partial_{x}, \quad b_{*}\left(s \partial_{u_{i}}\right)=x \partial_{y_{i}} .
$$

Thus if $A \in \Psi_{0}^{m, a, b}$, and $\kappa(A)=$ $k(s, u, \tilde{x}, \tilde{y})\left|\frac{d s d u d \tilde{x} d \tilde{y}}{s^{n^{\tilde{x}}}{ }^{n}}\right|^{1 / 2}=k \cdot r \quad$ then

$$
\kappa\left(x \partial_{x} \cdot A\right)=\left[\frac{\partial k}{\partial s}-\frac{n_{2}}{2}\right] \cdot \gamma
$$

(2.39)

$$
\kappa\left(x \partial_{y_{i}} \cdot A\right)=s \frac{\partial k}{\partial u_{i}} \cdot \gamma
$$

These new kernels each belong to $K_{0}^{m+1, a, b}$, so we are done.

A symbol map on Diff ${ }_{0}(M)$ was introduced in (2.7); there is a natural generalization of this map to $\Psi_{0}^{*, a, b}$. Recall first that

$$
\Psi_{0}^{m} \cap \Psi_{0}^{-\infty, a, b}=\Psi_{0}^{-\infty}
$$

This implies that if $A=A^{\prime}+A^{\prime \prime}$ is a decomposition as in
(2.35) ii), then it makes good sense to let ${ }^{\circ} \sigma_{m}(A)=$ ${ }^{o} \sigma_{m}\left(A^{\prime}\right)$ and only define the symbol on $\Psi_{0}^{m}$. Now, the map here, apart from simplifications in the density factors, is just the usual symbol on the space of distributions conormal along $\Delta \iota_{0}$ :

$$
\sigma_{\mathrm{m}}: \mathrm{K}_{0}^{\mathrm{m}}\left(\mathrm{M} ; \Gamma_{0}^{1 / 2}\right) \longrightarrow \mathrm{s}^{\mathrm{m}}\left(\mathrm{~N}^{*}\left(\Delta \iota_{0}\right) ; \Gamma_{0}(\mathrm{M}) \otimes \Gamma(\mathrm{fibre})\right)
$$

The space on the right contains symbolic densities on the fibres of the conormal bundle of $\Delta_{0} ; \quad \Gamma$ (fibre) -densities on the fibres of this bundle-arises through the use of the invariant Fourier transform. Also $\Gamma_{0}(M)$ results from the restriction of $\Gamma_{0}^{1 / 2}\left(\mathrm{M} \times{ }_{0} \mathrm{M}\right)$ to the diagonal, for the half densities on each factor combine along $\Delta c_{0}$ to give densities of weight one.

The $V_{0}$ symbol map ${ }^{0} \sigma_{m}$ is obtained by first applying the map $\sigma_{m}$ above, and then removing the density factors. It suffices to construct a natural density on $N^{*}\left(\Delta \iota_{0}\right)$ by which to 'divide'. This proceeds in three steps. First note the isomorphism

$$
\Gamma_{0}(\mathrm{M}) \otimes \Gamma(\text { fibre }) \simeq \Gamma_{0}\left(\mathrm{~N}^{*}\left(\Delta \iota_{0}\right)\right),
$$

the latter being the space of densities along $N^{*}\left(\Delta \iota_{0}\right)$ singular like $R^{-n}$ at $F$. Next, observe that there is a natural isomorphism

$$
\delta: N^{*}\left(\Delta \iota_{0}\right) \xrightarrow{\sim}{ }^{0} \mathrm{~T}^{*} \mathrm{M} .
$$

This is dual to the isomorphism $N\left(\Delta \iota_{O}\right) \xrightarrow{\sim}{ }^{0} \mathrm{TM}$, which is the usual natural isomorphism $N(\Delta \iota) \longrightarrow T M$ away from $\partial \Delta \iota_{0}$; near this boundary it comes from the identification of $x \partial_{x}, \quad x \partial_{y_{i}}$, which span ${ }^{0} T M$ near $\partial M$, with their lifts $s \partial_{s}, \quad s \partial_{u_{i}}$, which span $N\left(\Delta \iota_{O}\right)$ near $\partial \Delta \iota_{O}$. Finally, by way of the natural map

$$
\mathrm{T}^{*} \mathrm{M} \longrightarrow{ }^{\mathbf{0}} \mathrm{T}^{*} \mathrm{M}
$$

(which is not an isomorphism at $\partial M$ ) the symplectic density on $T^{*} M$ corresponds to a density on ${ }^{\circ} T^{*} M$. This density is singular at $\partial M$, yet a nonsingular section of $\Gamma_{0}\left({ }^{0} \mathrm{~T}^{*} \mathrm{M}\right)$. The resulting naturally defined element of $\Gamma_{0}(M) \otimes \Gamma(f i b r e)$ may thusly be removed so as to obtain the $\mathrm{V}_{\mathrm{O}}$ symbol map

$$
\begin{equation*}
{ }^{o^{\sigma_{m}}}: \Psi_{0}^{\mathrm{m}} \longrightarrow \mathrm{~S}^{\mathrm{m}}\left({ }^{\mathrm{o}} \mathrm{~T}^{*} \mathrm{M}\right) \tag{2.40}
\end{equation*}
$$

Here we have also used $\delta$ to transfer the symbols from $N^{*}\left(\Delta \iota_{0}\right)$ to ${ }^{0} T^{*} M$.
(2.41) Theorem: $\Psi_{0}^{*}$ is filtered by the symbol maps ${ }^{0}{ }_{\sigma_{m}}$, $m \in \mathbb{R}$. For any such $m$ we have a short exact sequence

$$
0 \longrightarrow \Psi_{0}^{\mathrm{m}-1} \longrightarrow \Psi_{0}^{\mathrm{m}} \xrightarrow{\mathrm{o}_{\mathrm{m}}} \mathrm{~S}^{\mathrm{m}}\left({ }^{\mathrm{o}^{*}} \mathrm{M}^{*}\right) / \mathrm{S}^{\mathrm{m}-1}\left({ }^{\mathrm{o}} \mathrm{~T}^{*} \mathrm{M}\right) \longrightarrow 0
$$

and the product formula

$$
{ }^{\circ} \sigma_{m_{1}+m_{2}}\left(A_{1} \cdot A_{2}\right) \equiv{ }^{o} \sigma_{m_{1}}\left(A_{1}\right)^{o} \sigma_{m_{2}}\left(A_{2}\right)
$$

Proof: The exactness of the sequence is a consequence of the same property of the symbol map on distributions conormal to $\Delta \iota_{0}$. The product formula follows from the usual one on standard pseudodifferential operators, provided we ascertain that $A_{1} \cdot A_{2}$ is in fact a $V_{0}$ pseudodifferential operator. For, both the symbol map and the composition rule agree with their ordinary analogues away from $\partial M$, hence the product rule must extend by continuity, provided the left hand side makes sense.

The first step is to derive an expression for the kernel of a composition. Use the representation (2.30), but write the kernels $k_{1}$, $k_{2}$ of $A_{1}$, $A_{2}$ as functions of $x, y, s, u$ (rather than $\tilde{x}, \tilde{y}, s, u)$. This is legitimate since we may assume that $k_{1}$ and $k_{2}$ are supported quite near $\Delta \iota_{0}$ With

$$
s_{1}=\frac{x}{\sim}, \quad u_{1}=\frac{\tilde{y}-\tilde{y}}{\tilde{x}}, \quad s_{2}=\frac{\tilde{x}}{\tilde{x}}, \quad u_{2}=\frac{\tilde{y}-\tilde{y},}{\tilde{x}}
$$

and recalling the half densities of (2.28), the composition $A_{1} A_{2} f$ becomes, on the level of kernels,

$$
\begin{aligned}
& f=f(\tilde{x} \cdot \tilde{y} \cdot) \cdot \mu \longmapsto \int k_{2}\left(\tilde{x}, \tilde{y}, s_{2}, u_{2}\right) f\left(\frac{\tilde{x}}{s_{2}}, \tilde{y}-\frac{\tilde{x}}{s_{2}} u_{2}\right) \frac{d s_{2}}{s_{2}} d u_{2} \cdot \mu \\
& \longmapsto \int k_{1}\left(x, y, s_{1}, u_{1}\right) k_{2}\left(\frac{x}{s_{1}}, y-\frac{x}{s_{1}} u_{1}, s_{2}, u_{2}\right) \\
& f\left(\frac{x}{s_{1} s_{2}}, y-\frac{x}{s_{1} s_{2}}\left(s_{2} u_{1}+u_{2}\right)\right) \cdot \frac{d s_{1} d u_{1} d s_{2} d u_{2}}{s_{1} s_{2}} \cdot \mu .
\end{aligned}
$$

Now, introduce

$$
\sigma=\frac{\mathrm{x}}{\tilde{x}^{\prime}}=\mathrm{s}_{1} \mathrm{~s}_{2}, \quad v=\frac{\frac{\mathrm{y}-\tilde{y^{\prime}}}{\tilde{x^{\prime}}}}{\mathrm{x}^{\prime}}=\mathrm{s}_{2} \mathrm{u}_{1}+\mathrm{u}_{2} .
$$

Hence

$$
\mathrm{A}_{1} \mathrm{~A}_{2} \mathrm{f}=\int \mathrm{k}(\mathrm{x}, \mathrm{y}, \sigma, v) \mathrm{f}\left(\frac{\mathrm{x}}{\sigma}, \mathrm{y}-\mathrm{x} v / \sigma\right) \frac{\mathrm{d} \sigma}{\sigma} \mathrm{~d} v \cdot \mu
$$

where, by the substitution $s_{1}=\sigma / s_{2}, \quad u_{1}=\left(v-u_{2}\right) / s_{2}$

$$
\begin{gathered}
\mathrm{k}(\mathrm{x}, \mathrm{y}, \sigma, v)= \\
\int \mathrm{k}_{1}\left(\mathrm{x}, \mathrm{y}, \sigma / \mathrm{s}_{2},\left(v-\mathrm{u}_{2}\right) / \mathrm{s}_{2}\right) \mathrm{k}_{2}\left(\sigma \mathrm{x} / \mathrm{s}_{2}, \mathrm{y}-\mathrm{x}\left(v-\mathrm{u}_{2}\right) / \sigma, \mathrm{s}_{2}, \mathrm{u}_{2}\right) \frac{\mathrm{ds} \mathrm{~s}_{2} \mathrm{~d} \mathrm{u}_{2}}{\mathrm{~s}_{2}^{\mathrm{n}-1}} .
\end{gathered}
$$

The $\mathbf{k}_{\mathrm{i}}$ are conormal along $\left\{\mathbf{s}_{\mathbf{i}}=1, \mathbf{u}_{\mathrm{i}}=0\right\}, \quad \mathbf{i}=1,2$, so we must show that $k$, as given above, is conormal along, and supported near, $\{\sigma=1, v=0\}$. The support condition is obvious from the similar assumption on $k_{1}$, $k_{2}$. As for the conormality, consider first the special case $x=y=0$ :
(2.42) Lemma: If $k_{1}, k_{2}$ are distributions on the group $G_{p}$ of (2.4) which are conormal at, and supported near, the identity (1,0). Then

$$
\mathrm{k}(\sigma, v)=\int \mathrm{k}_{1}\left((\sigma, v) \cdot(\mathrm{s}, \mathrm{u})^{-1}\right) \mathrm{k}_{2}(\mathrm{~s}, \mathrm{u}) \mathrm{s}^{-\mathrm{n}} \mathrm{~d} s \mathrm{du}
$$

enjoys the same properties.

Proof: We reduce this to the corresponding assertion for the group $\mathbb{R}^{\text {n }}$ where this result is well known. The main point is that the space of distributions conormal to some submanifold (e.g. ( 1,0 )) is coordinate invariant, and furthermore that, under a submersion, conormal distributions pull back to conormal distributions. Thus, since

$$
\begin{gathered}
(\sigma, v, \sigma-\mathrm{s}, v-\mathrm{u}) \longmapsto(\sigma / \mathrm{s},(v-\mathrm{u}) / \mathrm{s}) \\
\sigma / \mathrm{s}=\left(1-\frac{\sigma-\mathrm{s}}{\sigma}\right)^{-1}, \frac{v-\mathrm{u}}{\mathrm{~s}}=\frac{v-\mathrm{u}}{\sigma}\left(1-\frac{\sigma-\mathrm{s}}{\sigma}\right)^{-1}
\end{gathered}
$$

is a surjection, there is a distribution $k_{1}^{\prime}$ defined by

$$
\mathrm{k}_{1}\left(\sigma / \mathrm{s}, \frac{(v-\mathrm{u})}{\mathrm{s}}\right)=\mathrm{k}_{1}^{\prime}(\sigma, v, \sigma-\mathrm{s}, v-\mathrm{u})
$$

which is conormal at $\sigma-s=0, v-u=0$, and $C^{\infty}$ elsewhere. The pairing

$$
\int \mathrm{k}_{1}(\sigma, v, \sigma-\mathrm{s}, v-\mathrm{u}) \mathrm{k}_{2}(\mathrm{~s}, \mathrm{u}) \mathrm{s}^{-\mathrm{n}} \mathrm{~d} \mathrm{~s} \mathrm{du}
$$

defining $k$ is therefore of the required regularity, which proves the lemma.

To understand the group convolution when the parameters $x$, $y$ are present, only a slight modification is necessary. Thus, as before, there is a distribution

$$
\mathrm{k}_{2}^{\prime}\left(\mathrm{x}, \mathrm{y}, \sigma, v, \mathrm{~s}_{2}, \mathrm{u}_{2}\right)=\mathrm{k}_{2}\left(\mathrm{x}\left(1-\frac{\sigma-\mathrm{s}_{2}}{\sigma}\right)^{-1}, \mathrm{y}-\frac{\mathrm{x}}{\sigma}\left(v-\mathrm{u}_{2}\right), \mathrm{s}_{2}, \mathrm{u}_{2}\right)
$$

also conormal along $s_{2}=1, u_{2}=0$, and for which $x, y$, $\sigma, v$ are smooth parameters. As in the lemma, the kernel $\mathrm{k}(\mathrm{x}, \mathrm{y}, \sigma, v)$ is conormal along $\sigma=1, \quad v=0$. The proof of (2.41) is complete.

Remarks: Though we did not verify it, that ${ }^{0_{0} \sigma_{m}(P), ~ P \epsilon ~}$ Diff $f_{0}^{m}$, as defined here agrees with the earlier definition
(2.7) is quite easy to check. In addition, the product rule

$$
{ }^{o} \sigma_{m+\mu}(\mathrm{P} \cdot \mathrm{~A}) \equiv{ }^{\mathrm{o}} \sigma_{m}(\mathrm{P})^{o} \sigma_{\mu}(\mathrm{A}), \quad A \in \Psi_{0}^{\mu}
$$

is simpler to prove, since we may operate directly on the kernel of $A$ with the lifted operators on $M \times{ }_{0} M$ covering $P$.

The residual space for this symbol calculus is $\Psi_{0}^{-\infty}$; the kernels of the operators in this space are $C^{\infty}$ on the stretched product $M \times{ }_{0} M$, but not on $M \times M$. Thus, although smoothing in the interior of $M$, they are not compact on any reasonable space. This is discussed at length in the next two sections.

## D. Indicial and Normal Operators

The first step in constructing a parametrix for $P \in$ Diffo is to iteratively 'remove' the conormal singularities of its kernel along $\Delta \iota_{0}$. As in the ordinary construction, this is reduced to algebraic manipulation by way of the symbol calculus. But, as noted at the end of the last section, the resulting error after this step is not compact. This necessitates the use of two additional steps, both iterative in nature, and both employing models for the various operators involved. These we describe now.

If $E_{O}$ is the operator obtained in the first step:

$$
P E_{O}=I-Q_{O}, \quad Q_{O} \in \Psi_{0}^{-\infty}
$$

then $\kappa_{0}$, the kernel of $Q_{0}$, is $C^{\infty}$ on $M \times_{O} M$ up to all boundaries. In fact, $\kappa_{0}$ vanishes to infinite order at $T$ and $B$; at $F$, however, its Taylor series is nontrivial (which causes $Q_{0}$ not to be compact). Thus the second step of the construction is the removal of this Taylor series. We seek an $E_{1}$ such that

$$
P E_{1}=Q_{0}-Q_{1}
$$

and $\kappa_{1}$, the kernel of $Q_{1}$, vanishes to infinite order at $F$. Unfortunately, in order to accomplish this, we must settle for a correction term $E_{1}$ lying in $\Psi_{0}^{-\infty, a, b}$ for some $a, b$ (We are restricting discussion to scalar operators for simplicity.) Consequently

$$
\mathrm{Q}_{1} \in \mathrm{R}^{\infty} \Psi_{0}^{-\infty, a, b}, \quad \kappa_{1} \in R^{\infty} K_{0}^{-\infty, a, b}
$$

where the space of kernels on the right consists of those elements of $K_{0}^{-\infty}, a, b$ which vanish to infinite order on $F$. If $a$ and $b$ are sufficiently large then $Q_{1}$ is actually
compact and the construction is complete. Usually, though, it is necessary to carry through the third step: the removal of the conormal singularity of $\kappa_{1}$ along the top face. Here we seek an $E_{2} \in R^{\infty} \Psi_{0}^{-\infty}, a, b$ such that

$$
P E_{2}=Q_{1}-Q_{2}, \quad Q_{2} \in R^{\infty} \Psi_{0}^{-\infty, \infty, b}
$$

The operator $Q_{2}$ is compact on appropriate weighted $L^{2}$ spaces and

$$
P \cdot\left(E_{0}+E_{1}+E_{2}\right)=I-Q_{2}
$$

so that $E=E_{0}+E_{1}+E_{2}$ is the desired parametrix.
From this somewhat vague discussion it should be clear that $E_{1}$ and $E_{2}$ may be 'locally' constructed. Specifically, $E_{1}$ should only depend on the operator $P$ frozen at $F$ and the relevant Taylor series of $\kappa_{0}$, whereas $E_{2}$ must rely only on the infinitesimal action of $P$ near $T$ and the conormal expansion of $\kappa_{1}$ there.

Consider first the former of these steps. Toward the stated goal, define the normal operator of $A \in \Psi_{0}^{-\infty, a, b}$

$$
\begin{equation*}
N_{p}(A)=\left.\kappa(A)\right|_{F_{p}} \tag{2.43}
\end{equation*}
$$

By the remark following (2.26), and recalling that $k(A)$ is actually a section of the half density bundle, $N_{p}(A)$
is a section of

$$
\Gamma_{0}^{1 / 2}\left(F_{p}\right) \otimes \not A_{p h g}^{a, b}\left(F_{p}\right)
$$

i.e. smooth in the interior, and with the appropriate conormal singularity at $T \cap F_{p}$ and $B \cap F_{p}$. As it stands, this seems to bear no relationship to the previous definition (2.8) for $N_{p}(P), \quad P \in \operatorname{Diff}_{0}$. To make this comparison we must interpret $N_{p}(A)$ as an operator. Remember that $F_{p}$ is naturally identified (in its interior) with the group $G_{p}$ in two ways. Using the $G_{p}^{\ell}$ identification, $N_{p}(A)$ is to be thought of as a left convolution operators on half densities over $\mathrm{F}_{\mathrm{p}}$. Thus, in the coordinates $(s, u, \tilde{x}, \tilde{y})$ of $(2.20), \quad p=(0,0):$

$$
N_{p}(A)=k(s, u, 0,0) \cdot\left|\frac{d s d u d \tilde{s} d \tilde{u}}{s^{n} \mathbf{s}^{n}}\right|^{1 / 2}
$$

(2.44)

$$
N_{p}(A) f=\int k\left((s, u) \cdot(\tilde{s}, \tilde{u})^{-1}, 0,0\right) f(\tilde{s}, \tilde{u}) \frac{d \tilde{s} d \tilde{u}}{\tilde{s}^{n}} \cdot\left|\frac{d s d u}{s^{n}}\right|^{1 / 2}
$$

where $f=f(s, u)\left|d s d u / s^{n}\right|^{1 / 2}$.
Now, if $P \in D_{i f f}^{m}$ is expressed as in (2.2), then by (2.21) its lift to $M \times{ }_{0} M$ is
a) $\sum_{|\alpha| \leq m} a_{\alpha}(s \tilde{x}, \tilde{y}+\tilde{x} u)\left(s \partial_{u}\right)^{\alpha}\left(s \partial_{s}\right)^{\alpha}{ }^{n}, \quad \alpha=\left(\alpha^{\prime}, \alpha_{n}\right)$.

Hence, applied to the kernel of the identity (2.31),

$$
\kappa(P)=\sum_{|\alpha| \leq m} a_{\alpha}(s \tilde{x}, \tilde{y}+\tilde{x} u)\left(s \partial_{u}\right)^{\alpha}\left(s \partial_{s}-\frac{n}{2}\right)^{\alpha} n_{\delta(s-1) \delta(u) \cdot \gamma}
$$

so that, according to (2.43)
b) $N_{p}(P)=$

$$
\sum_{|\alpha| \leq m} a_{\alpha}(0,0)\left(s \partial_{u}\right)^{\alpha}\left(s \partial_{s}-\frac{n}{2}\right)^{\alpha} n_{\delta(s-1) \delta(u)}\left|\frac{d s d u d \tilde{s} d \tilde{u}}{s^{n^{2} \tilde{s}^{n}}}\right|^{1 / 2}
$$

However, according to (2.8)

$$
\text { c) } N_{p}(P)=\sum_{|\alpha| \leq m} a_{\alpha}(0,0)\left(s \partial_{u}\right)^{\alpha}\left(s \partial_{s}\right)^{\alpha}{ }^{n} .
$$

Finally then it is not hard to see that the expression $c$ ). applied to a half density $f(s, u)\left|d s d u / s^{n}\right|^{1 / 2}$, coincides with the result of setting expression b) into (2.44). We have shown that (2.43) is a proper generalization of (2.8).
(2.45) Proposition: There is an exact sequence

$$
0 \longrightarrow R \Psi_{0}^{-\infty, a, b} \longrightarrow \Psi_{0}^{-\infty, a, b} \xrightarrow{N_{p}}
$$

$$
\alpha_{\mathrm{phg}}^{\mathrm{a}, \mathrm{~b}}\left(\mathrm{~F}_{\mathrm{p}}\right) \otimes \Gamma_{0}^{1 / 2}\left(\mathrm{~F}_{\mathrm{p}}\right) \longrightarrow 0 .
$$

If $P \in \operatorname{Diff}_{0}^{m}, \quad A \in \Psi_{0}^{-\infty}, a, b$ then

$$
N_{p}(P \cdot A)=N_{p}(P) \cdot N_{p}(A)
$$

Remark: For $k=0,1,2, \cdots, \quad R^{k_{\Psi}} \Psi_{0}^{-\infty}, a, b$ is the space of operators whose kernels vanish to order $k$ at $F$.

Proof: The exactness of the sequence is obvious, save perhap the surjectivity of $N_{p}$. But that too is easy since one may readily construct in local coordinates a kernel with fixed restriction to $F_{p}$. As for the product formula, either by applying the operator a) above to $\kappa(A)$ and then setting $\tilde{x}=\tilde{y}=0$, or by applying the operator $c$ ) to $N_{p}(A) \quad$ in $(2.44)$ one arrives at the same section of $\Gamma_{0}^{1 / 2}$, for $\tilde{x}$ and $\tilde{y}$ are merely parameters in a).

The conjugate $R^{k} N_{p} R^{-k}$ acts on $R^{k_{4}}{ }_{0}^{-\infty, a, b}$, though we shall not use this fact. Nonetheless, as $k$ varies in $\mathbb{Z}^{+}$, the residual space

$$
R^{\infty} \Psi_{0}^{-\infty}, a, b
$$

is the outer limit of the effectiveness of $N_{p}$. Kernels of elements in this space, when pushed down to $M \times M$, lie in

$$
\alpha_{0}^{a, b}(M \times M) \otimes \Gamma_{0}^{1 / 2}(M \times M)
$$

The space $\mathscr{A}_{0}^{a, b}$ consists of those distributions which are conormal (with expansions) along the boundary components $\partial M \times M, M \times \partial M$ of $M \times M$, and which vanish to all orders at $\partial \Delta \iota$.

The final step utilizes the indicial operator, which once again models kernels in the last residual space. It is defined by the exact sequence

$$
\begin{align*}
0 \longrightarrow R^{\infty} \Psi_{0}^{-\infty}, a+1, b & R^{\infty} \Psi_{0}^{-\infty, a, b} \xrightarrow{I}  \tag{2.46}\\
& A_{0}^{a, b} \otimes \Gamma_{0}^{1 / 2} / A_{0}^{a+1, b} \otimes \Gamma_{0}^{1 / 2} \longrightarrow 0
\end{align*}
$$

Thus, for $A \in R^{\infty} \Psi_{0}^{-\infty, a, b}$, $I(A)$ is simply the top term of the expansion of the kernel of $A$ on $\partial M \times M$. It too may be defined as an operator, but we have no need of this.

On the other hand, $P \in \operatorname{Diff}_{0}^{m}(M)$ also has an indicial operator, which we do want to regard as an operator. Note that

$$
A \in R^{\infty} \Psi_{0}^{-\infty}, a, b \Rightarrow P \cdot A \in R^{\infty} \Psi_{0}^{-\infty}, a, b
$$

I(P) is now defined to be that operator acting
‘infinitesimally at $\partial M$ ’, or equivalently on $N \partial M$, for which

$$
\begin{equation*}
I(P \cdot A)=I(P) I(A) \tag{2.47}
\end{equation*}
$$

Strictly speaking there are many operators satisfying (2.47). But since

$$
\begin{aligned}
& x^{\mu} \partial_{x}: \mathscr{A}_{0}^{\mathrm{a}, \mathrm{~b}} \longrightarrow \mathscr{A}_{0}^{\mathrm{a}+1, \mathrm{~b}} \quad \mu>1 \\
& \mathrm{x}^{v} \partial_{\mathrm{y}_{\mathrm{i}}}: \mathscr{A}_{0}^{\mathrm{a}, \mathrm{~b}} \longrightarrow \mathscr{A}_{0}^{\mathrm{a}+1, \mathrm{~b}}
\end{aligned}
$$

there is a natural choice for $I(P)$. If

$$
P=\sum_{k+|\alpha| \leq m} a_{k, \alpha}(x, y)\left(x \partial_{x}\right)^{k}\left(x \partial_{y}\right)^{\alpha}
$$

then set

$$
\begin{equation*}
I(P)=\sum_{k \leq m} a_{k, 0}(0, y)\left(x \partial_{x}\right)^{k} \tag{2.48}
\end{equation*}
$$

This is well defined, provided only coordinate systems of the type discussed in section $C$ are used, and obviously satisfies (2.47). The residual space of the filtration induced by $I$ is

$$
\mathrm{R}^{\infty} \Psi_{0}^{-\infty, \infty, b}
$$

the kernels of which are compact on $L^{2}$ (so long as $b$ is large enough), as we prove in the next section.

## E. Continuity Properties

We now study the continuity properties of $V_{0}$ pseudodifferential operators. Following the usual scheme of such proofs, first the $L^{2}$ boundedness of residual operators is demonstrated-this is the most arduous part of the proofthen the symbol calculus immediately implies boundedness for operators of order zero. Finally a few more refined statements, including criteria for compactness, are proved. Other less general results are discussed in the last section of the next chapter. As before, we limit all consideration to scalar operators; the analogues for operators on the form bundles are either obvious or will be commented on when required.
(2.49) Proposition: Any $A \in \Psi_{0}^{-\infty, a, b}$ induces a bounded map

$$
A: x^{r_{L}^{2}}\left(x^{-n} d x d y\right) \longmapsto x^{r^{\prime}} L_{l o c}^{2}\left(x^{-n} d x d y\right)
$$

provided $r^{\prime} \leq r, a-r>\frac{n-1}{2}, b+r^{\prime}>\frac{n-1}{2}, \quad$ and $a+b>n-1$.

Proof: The slightly bizarre method we employ is to reduce the statement to the familiar one concerning $L^{2}$ boundedness of ordinary pseudodifferential operators of order zero. First note that we may assume that the kernel k of $A$ is supported in $0 \leq x, \tilde{x} \leq 1,|y|,|\tilde{y}| \leq 1$. It is of course $C^{\infty}$ when lifted to $M \times{ }_{0} M$, except at $T$ and B where it is conormal. It also suffices to demonstrate boundedness of

$$
x^{-r \cdot-(n-1) / 2} \underset{A x^{r}+(n-1) / 2}{\sim} \quad \text { on } \quad L_{c}^{2}\left(x^{-1} d x d y\right)
$$

The kernel of this last operator is

$$
k_{c}=x^{r-r^{\prime}}\left[\frac{x}{\sim}\right]^{-r-(n-1) / 2} k=\tilde{x}^{r-r} \cdot\left[\frac{\tilde{x}}{x}\right] r^{\prime}+(n-1) / 2 k
$$

Split $k_{c}$ into a sum of two kernels, one supported in $x \leq$ $2 \tilde{x}$ and the other in $\tilde{x} \leq 2 x$. As $r-r \geq 0$, let $k_{1}$ be the first of these divided by $x^{r-r^{\prime}}$ and $k_{2}$ the second divided by $\tilde{x}^{r-r^{\prime}}$. We need only demonstrate boundedness of $\mathrm{k}_{1}$ and $\mathrm{k}_{2}$ since $\mathrm{k}_{1}+\mathrm{k}_{2}$ dominates $\mathrm{k}_{\mathrm{c}}$.

$$
\text { Write } k_{1} \text { as a function of } \tilde{x}, y, s, u \text { as in (2.30) }
$$

and $k_{2}$ as a function of $x, y, t, v$ as in (2.32). Then by assumption $k_{1}$ is supported in $0 \leq \tilde{x} \leq 1, \quad|y| \leq 1$,
$0 \leq s \leq 2, \quad$ and $k_{2}$ in $0 \leq x \leq 1, \quad|y| \leq 1,0 \leq t \leq 2 ;$ both are supported within the domain of regularity of their coordinate systems. Furthermore

> a) $\mathrm{k}_{1} \in \mathrm{C}^{\infty}$ for $\mathrm{s} \neq 0$, is conormal at $s=0$ with an expansion of lowest power $s^{a-r-(n-1) / 2}$ and a symbol of order $-a-b$ in $u$
b) $k_{2} \in C^{\infty}$ for $t \neq 0$, is conormal at $t=0$ with an expansion of lowest power $t^{b+r^{\prime}+(n-1) / 2}$ and a symbol of order - $\mathrm{a}-\mathrm{b}$ in u .

Consider first $k_{1}$ alone. By the Plancherel formula

$$
\begin{aligned}
& \int k_{1}\left[\tilde{x}, y, \frac{x}{\sim}, \frac{y-\tilde{y}}{\tilde{x}}\right] f(\tilde{x}, \tilde{y}) \tilde{x}^{-n} d \tilde{x} d \tilde{y}= \\
& \int \hat{\mathrm{k}}_{1}[\tilde{\mathrm{x}}, \mathrm{y}, \underset{\tilde{x}}{\underset{\mathrm{x}}{x}},-\tilde{\mathrm{x}} \eta] \mathrm{e}^{-\mathrm{i} \mathrm{y} \cdot \eta} \hat{\mathbf{f}}(\tilde{\mathrm{x}}, \eta) \tilde{\mathrm{x}}^{-1} \mathrm{~d} \tilde{\mathbf{x}} \mathrm{~d} \eta
\end{aligned}
$$

since the Fourier transform of $k_{1}\left[-\frac{y-\tilde{y}}{\tilde{x}}\right]$ with respect to $\tilde{y}$ is $\hat{\mathrm{k}}_{1}(\tilde{x}, \mathrm{y}, \mathrm{x} / \tilde{\mathrm{x}},-\mathrm{x} \eta) \mathrm{e}^{-\mathrm{i} y \cdot \eta_{\mathrm{x}} \tilde{n}^{n-1}}$. Next, introduce into this last integral the new variables $e^{-t}=x, e^{-\tilde{t}}=\tilde{x}$ to obtain

$$
\int \hat{k}_{1}\left(e^{-\tilde{t}}, y, e^{-(t-\tilde{t})},-e^{-\tilde{t}} \eta\right) e^{-i y \cdot \eta} \hat{f}\left(e^{-\tilde{t}}, \eta\right) d \tilde{t} d \eta
$$

where the integrand is supported in $t \geq 0, \tilde{t} \geq 0$, $\mathrm{t}-\tilde{\mathrm{t}} \geq-1$. Notice that by (2.50) a) the first factor increases exponentially unless $a-r>(n-1) / 2$, in which case it decreases rapidly. Finally set

$$
\kappa_{1}(\tilde{t}, y, \tau, \eta)=\int \hat{\mathrm{k}}_{1}\left(\mathrm{e}^{-\tilde{\mathrm{t}}}, \mathrm{y}, \mathrm{e}^{-\sigma}, \mathrm{e}^{-\tilde{\mathrm{t}}} \eta\right) \mathrm{e}^{-\mathrm{iy} \cdot \eta} \mathrm{e}^{-\mathrm{i} \sigma \tau} \mathrm{~d} \sigma
$$

so that the former integral becomes

We shall show that $\kappa_{1}$ is a symbol in $\tau, \eta$ of order 0 (and type (1,0)), except for a mild conormal singularity at $e^{-\tilde{t}} \eta=0$, and so represents an $L^{2}(d t d y)$ bounded operator. Since $L^{2}(d t d y)=L^{2}\left(x^{-1} d x d y\right)$, this shows that $k_{1}$ has the required boundedness.

To prove the claim we need only examine the effect of differentiating the integral which defines $\kappa_{1}$. First notice that by (2.50) a), the symbolic properties of $k_{1}(\tilde{x}, y, s, u)$ in $u$ imply that its Fourier transform is rapidly decreasing in the dual variable $\zeta$ and conormal at $\zeta=0$. Since $a+b>(n-1), \hat{k}_{1}$ is bounded there. Hence $\kappa_{1}$ also enjoys these properties. In addition, if a-r > ( $n-1$ )/2 as assumed, $k_{1}$ is rapidly decreasing in $T$. Thus differentiating with respect to either y or $\boldsymbol{T}$
preserves this rapid decrease. Differentiating with respect to $\tilde{t}$ either affects the first slot of $\hat{k}_{1}$ or the fourth. As $\tilde{x}$ is a smooth parameter, the former preserves regularity. In the later case a factor $e^{-\tilde{t}} \eta$ is produced in front:

$$
\partial_{\mathrm{t}^{\kappa}}^{\sim}=-\int\left(\mathrm{e}^{-\tilde{\mathrm{t}}} \partial_{\tilde{x}} \hat{\mathrm{k}}_{1}+\mathrm{e}^{-\tilde{\mathrm{t}}} \eta_{\mathrm{j}} \partial_{\mathrm{j}} \hat{\mathrm{k}}_{1}\right) \mathrm{e}^{-\mathrm{iy} \cdot \eta-\mathrm{i} \sigma \tau} \mathrm{~d} \sigma
$$

is of the same regularity as $\kappa_{1}$, in particular bounded. Finally,

$$
\eta_{i} \partial_{\eta_{j}} \kappa_{1}=\int-\left(\eta_{i} e^{-\tilde{t}^{2}} \partial_{\zeta_{j}} \hat{\mathrm{k}}_{1}+\mathrm{i} \eta_{\mathrm{i}} \mathrm{y}_{\mathrm{j}} \hat{\mathrm{k}}_{1}\right) \mathrm{e}^{-\mathrm{iy} \cdot \eta-\mathrm{i} \sigma \tau} \mathrm{~d} \sigma
$$

is also uniformly bounded. By iteration it is readily seen that

$$
\left|\partial_{\mathrm{t}}^{\mathrm{j}} \partial_{\mathrm{y}}^{\alpha} \partial_{\tau}^{\ell} \partial_{\eta}^{\beta}{ }_{\kappa_{1}}(\tilde{\mathrm{t}}, \mathrm{y}, \tau, \eta)\right| \leq \mathrm{C}_{\mathrm{j} \mathrm{\ell} \mathrm{\alpha} \mathrm{\beta}}(1+|\tau|+|\eta|)^{-\ell-|\beta|}
$$

and, as mentioned before, $k_{1}$ has a bounded conormal singularity at $\eta=0$, uniform in all other parameters. By the standard theory $k_{1}$ induces a bounded transforma$t$ ion on $L^{2}\left(x^{-1} d x d y\right)$.

The proof for $k_{2}$ is quite similar, only here one requires the hypothesis $b+r^{\prime}>\frac{n-1}{2}$. This completes the proof of (2.49).

Before proving continuity of operators of order zero, let us note that $\Psi_{0}^{m}$ is invariant under the operation of taking adjoints, for any $m$. Thus if $A \in \Psi_{0}^{m}$ then the kernel of $A^{*}$ is

$$
\begin{aligned}
& \kappa_{A^{*}}(x, y, \tilde{x}, \tilde{y})\left|\frac{d x d y d \tilde{x} d \tilde{y}}{x^{n} \tilde{x}^{\tilde{n}}}\right|^{1 / 2}= \\
& \bar{\kappa}_{A}(\tilde{x}, \tilde{y}, x, y)\left|\frac{d x d y d \tilde{x} d \tilde{y}}{x^{n} \tilde{x}^{n}}\right|^{1 / 2}
\end{aligned}
$$

which certainly has the same regularity as $\kappa_{A}$ on $M x_{0} M$. In addition, it is easy to see that

$$
o_{\sigma_{m}}\left(A^{*}\right)=\overline{o_{\sigma(A)}}
$$

(2.51) Theorem: $A \in \Psi_{0}^{0, a, b}$ is bounded from

$$
\mathrm{x}^{\mathrm{r}} \mathrm{~L}^{2} \Gamma_{0}^{1 / 2} \longrightarrow \mathrm{x}^{\mathrm{r}^{\prime}} \mathrm{L}^{2} \Gamma_{0}^{1 / 2}
$$

provided $a+b>n-1, \quad a-r>\frac{n-1}{2}, \quad b+r^{\prime}>\frac{n-1}{2}, \quad r \geq r^{\prime}$.

Proof: Write $A=A_{1}+A_{2} \in \Psi_{0}^{0}+\Psi_{0}^{-\infty, a, b} . \quad A_{2}$ is bounded between these spaces by (2.49). As for $A_{1}$, use the symbol calculus of (2.41) to find $B \in \Psi_{0}^{O_{0}}$ such that for some sufficiently large $M \in \mathbb{R}$

$$
A_{1}^{*} A_{1}+B^{*} B=M \cdot I+R, \quad R \in \Psi_{0}^{-\infty} .
$$

Then, for $f=f(x, y)\left|x^{-n} d x d y\right|^{1 / 2}$

$$
\int\left|A_{1} f\right|^{2}+\int|B f|^{2}=M \int|f|^{2}+\int R f \cdot f
$$

which implies, using the boundedness of $R$, that

$$
\int\left|A_{1} f\right|^{2} \leq M^{\prime} \int|f|^{2}
$$

as desired.
The next step is to consider the action on Sobolevtype spaces. Since the interest here is with parametrices of differential operators, the full theory shall not be developed. For $m=0,1,2, \cdots$ set

$$
\begin{align*}
& \text { (2.52) } H_{b}^{m}\left(M, \Gamma_{0}^{1 / 2}\right)=\left\{u \in L^{2} \Gamma_{0}^{1 / 2}:\left(x \partial_{x}\right)^{i}\left(x \partial_{y}\right)^{\alpha} u \in L^{2},\right.  \tag{2.52}\\
& \\
& i+|\alpha| \leq m\} \\
& \text { (2.53) Corollary: } A \in \Psi^{-m, a, b} \text { is bounded from } \\
& \qquad x^{r} L^{2} \Gamma_{0}^{1 / 2} \longrightarrow x^{r^{\prime}} H_{b}^{m}\left(\Gamma_{0}^{1 / 2}\right) \\
& \text { for } r, r, a, b \text { as in }(2.51) .
\end{align*}
$$

Proof: If $k$ is the kernel of $A$, then recalling (2.39)

$$
\left(x \partial_{x}\right)^{i}\left(x \partial_{y}\right)^{\alpha} k \in K_{0}^{0, a, b} \quad \text { for } \quad i+|\alpha| \leq m
$$

The corollary follows.

The space of conormal half densities

$$
\begin{aligned}
& \mathscr{A}^{a}\left(M ; \Gamma_{0}^{1 / 2}\right)=\left\{f \in \Gamma_{0}^{1 / 2}: f=f(x, y)\left|\frac{d x d y}{x^{n}}\right|^{1 / 2},\right. \\
& \qquad\left(x \partial_{x}\right)^{i}\left(x \partial_{y}\right)^{\alpha} f(x, y) \in \rho^{a}(\log \rho)^{N_{L}}{ }^{\infty} \\
& \text { for all } i, \alpha, \text { and some fixed } N\}
\end{aligned}
$$

has, as in (2.34), the subspace of polyhomogeneous elements
(2.54) $\quad A_{\mathrm{phg}}^{\mathrm{a}}\left(\mathrm{M} ; \Gamma_{0}^{1 / 2}\right)=$

$$
\left\{f \in A^{a}: f \sim \sum_{i=0}^{\infty} \sum_{j=0}^{N_{i}} f_{i}(y) x^{a+i}(\log x)^{j}\right\}
$$

This is the natural range of $V_{0}$ pseudodifferential operators applied to smooth functions. Again we only prove a special case:
(2.55) Proposition: For any $b, m \in \mathbb{R}, A \in \Psi_{0}^{m, a, b}$ maps

$$
\dot{\mathrm{C}}^{\infty} \longrightarrow A_{\mathrm{phg}}^{\mathrm{a}}
$$

(where the space on the left denotes the smooth half densities vanishing to infinite order on $\partial M$ ).

Proof: If $k \in K_{0}^{m, a, b}$ is the kernel of $A$, then the product

$$
k(x, y, \tilde{x}, \tilde{y})\left|\frac{d x d y d \tilde{x} d \tilde{y}}{x^{n} \tilde{x}^{n}}\right|^{1 / 2} f(\tilde{x}, \tilde{y})\left|\frac{d \tilde{x} d \tilde{y}}{\tilde{x}^{n}}\right|^{1 / 2}
$$

$f \in \dot{C}^{\infty}$, vanishes to infinite order on $\{\tilde{x}=0\}$ and has an ordinary conormal singularity along the diagonal. Its integral then is that of an ordinary pseudodifferential kernel against f. Furthermore, if

$$
k \sim \sum k_{i}, \quad k_{i} \in{\underset{p h g}{a+i}, b}_{a}^{a}
$$

then each

$$
\int k_{i}(x, y, \tilde{x}, \tilde{y}) f(\tilde{x}, \tilde{y}) \frac{d \tilde{x} d \tilde{y}}{\tilde{x}^{n}}\left|\frac{d x d y}{x^{n}}\right|^{1 / 2} \epsilon A^{a+i}
$$

which gives the result.

Our final result concerns compactness of operators of negative order with appropriate boundary conditions.
(2.56) Proposition: $A \in R^{\infty} \Psi_{0}^{-m, a, b}$ is a compact map between

$$
\mathrm{x}^{\mathrm{r}} \mathrm{~L}^{2} \Gamma_{0}^{1 / 2} \longrightarrow \mathrm{x}^{\mathrm{r}^{\prime}{ }_{\mathrm{L}}{ }^{2} \Gamma_{0}^{1 / 2}}
$$

for $m>0, a+b>n-1, a-r>\frac{n-1}{2}, b+r^{\prime}>\frac{n-1}{2}$.

Proof: $\quad x^{-\epsilon} A=B \in \Psi_{0}^{-m, a-\epsilon, b}$, and if $\epsilon$ is sufficiently small, $B$ is bounded between

$$
x^{r^{\prime}}{ }^{2} \longrightarrow x^{r^{\prime}}{ }_{b}^{m}
$$

also. The compactness of $A=x^{\epsilon} B$ is furnished by the uniform smallness of $x^{\epsilon} B$ in a neighbourhood of $\partial M$ and the uniform equicontinuity of functions in $H_{b}^{m}$ on $\left\{(x, y): x \geq \epsilon_{1}\right\}$ for each $\epsilon_{1}>0$, cf. [29].

## A. Identification of the Model Operators

The novel step in constructing a parametrix for $\Delta_{g}$, as indicated in section $D$ of the previous chapter, is the use of simpler models for this operator at the boundary, where it degenerates. These models, the normal and indicial operators, are employed in iterative schemes analogous to the initial symbolic step of the construction to obtain a parametrix for which the remainder is compact. It is important then to reach a thorough understanding of their mapping properties; this is the goal of the present chapter. In this section we identify these operators explicitly and analyze the indicial operator. Since it is an ordinary differential operator of Euler type, this is quite easy. The normal operator, however, is a partial differential operator, and the derivation of those of its properties we need later will require more of an effort.

We have given two different, yet equivalent, definitions of the normal operator of $P \in \operatorname{Diff}_{0}^{m}(M)$. The first one, (2.8), is an operator on the half tangent space $M_{p}$, $p \in \partial M$, while the second, (2.43), is a convolution kernel on the fibre $F_{p}$ of the front face. This latter interpretation is the one we ultimately use, but the former is easier to manipulate, hence the one of interest in this chapter. Recall now that the action of $\Delta_{g}$ on ${ }^{o} \Lambda^{k}(M)$ is given by
the operator $P$ of (1.21), so that $P=\rho^{k} \Delta_{g} \rho^{-k}$. Then $N_{p}\left(\Delta_{g}\right)$ is an operator on $M_{p}$ acting on sections of ${ }^{o} \Lambda^{k}\left(M_{p}\right)$; the induced operator is $N_{p}(P)$. It turns out that $N_{p}\left(\Delta_{g}\right)$ is actually the Laplacian of a constant curbature metric, whence it has greater invariance properties than merely the expected $G_{p}$-invariance.
(3.1) Proposition: $N_{p}\left(\Delta_{g}\right)$ is the Laplacian of the constand curvature metric $\left(d \rho_{p}\right)^{-2} h_{p}$; here $d \rho_{p}$ is to be thought of as a linear function on $M_{p}$. The induced action on ${ }^{o} \Lambda^{k}\left(M_{p}\right)$ is given by $N_{p}(P)$.

Proof: Choose coordinates $\bar{z}=\left(\bar{z}^{1}, \cdots, \bar{z}^{n-1}, \bar{z}^{n}\right)=(\bar{y}, \bar{x})$ on $M$ near $p$, such that $\bar{y}=\left(\bar{z}^{1}, \cdots, \bar{z}^{n-1}\right)$ are geodesic coordinates for $h$ restricted to $\partial M$, centered at $p$, and $\bar{z}^{n}=\bar{x}$ is a defining function for the boundary and $\left|\partial / \partial \bar{z}^{\mathrm{n}}\right|_{\mathrm{h}}=1$. Let $\mathrm{z}=(\mathrm{y}, \mathrm{x})$ be the corresponding linear coordinates on $M_{p}$. Finally, write

$$
\rho(\bar{y}, \bar{x})=a(\bar{y}) \bar{x}+b(\bar{y}, \bar{x}), \quad b=O\left(\bar{x}^{2}\right) .
$$

Then $\left|d \rho_{p}\right|^{2}=a(0)^{2}$. Now, for $\omega \in \Omega^{k}(M)$

$$
\mathrm{P} \omega=\rho^{2} \Delta_{\mathrm{h}} \omega+(2-\mathrm{k}) \rho \mathrm{d} \rho \wedge \delta_{\mathrm{h}} \omega-\mathrm{k} \rho \delta_{\mathrm{h}}(\mathrm{~d} \rho \Lambda \omega)
$$

(1)
(2)
(3)

$$
\begin{equation*}
+(\mathrm{n}-\mathrm{k}) \rho \mathrm{L}_{\nabla \rho} \omega-2 \rho \iota \iota_{\nabla \rho} \mathrm{d} \omega \tag{5}
\end{equation*}
$$

(4)

$$
+(\mathrm{n}-2 \mathrm{k}) \mathrm{d} \rho \wedge \iota_{\nabla \rho}{ }^{\omega-\mathrm{k}(\mathrm{n}-\mathrm{k}-1) \iota_{\nabla \rho}(\mathrm{d} \rho) \omega}
$$

(6)
(7)
is the operator of (1.21), defined by $\Delta_{g}\left[\frac{\omega}{\rho k}\right]=\frac{(P \omega)}{\rho^{k}}$. Suppose $\quad \omega=\omega_{I} d^{-I}, \quad I=\left(i_{1}, \cdots, i_{k}\right), \quad 1 \leq i_{1}<\cdots<$ $i_{k} \leq n$. We shall consider each term (1) - (7) separately, always using the product rule (2.10).
(1) $\rho^{2} \Delta_{h} \omega$ has normal operator

$$
\left[-a(0)^{2} \sum_{i=1}^{n} \frac{\partial^{2} \omega_{I}}{\left(\partial z^{i}\right)^{2}}\right] d x^{I}
$$

(2) $\rho \mathrm{d} \rho \wedge \delta_{h} \omega=$

$$
\begin{aligned}
\left(a \bar{z}^{n}+b\right)\left(a d \bar{z}^{n}+\bar{z}^{n} d a\right. & +d b) \wedge\left(-\delta_{i_{s}}^{I}, I \backslash i_{s}\left(\nabla_{i} \omega_{I}\right) d \bar{z}^{I \backslash i_{s}}\right) \\
& =a^{2} \bar{z}^{n} d \bar{z}^{n} \wedge(-1)^{s} \frac{\partial \omega_{I}}{\partial z_{s}} d \bar{z}^{I \backslash i_{s}}+\ldots
\end{aligned}
$$

$$
-a^{2} z^{n} \frac{\partial \omega}{\partial z^{n}} \quad d z^{I} \quad \text { if } \quad i_{k}=n
$$

or

$$
\sum_{s} a^{2}(-1)^{k+s-1} z^{n} \frac{\partial \omega}{\partial z^{i}} d z^{I \backslash i} s d^{n} \quad \text { if } \quad i_{k}<n
$$

(3) $\rho \delta_{h}(d \rho \wedge \omega)=\left(a \bar{z}^{\mathrm{n}}+\mathrm{b}\right) \delta_{h}\left(\left(\operatorname{ad} \bar{z}^{\mathrm{n}}+\cdots\right) \wedge \omega_{\mathrm{I}} \mathrm{d}^{-\mathrm{z}}\right)$

$$
\begin{aligned}
= & \left(a \bar{z}^{-n}+b\right) \sum(-1)^{k+1} \delta_{i}^{I, n}, I \backslash i_{s}, n^{\nabla_{i}}\left(a \omega_{I}\right) d \bar{z}^{I \backslash i_{s}} \wedge d \bar{z}^{n} \\
& +\left(a \bar{z}^{n}+b\right)(-1)^{k+1} \delta_{n, I}^{I, n} \nabla_{n}\left(a \omega_{I}\right) d \bar{z}^{I}+\cdots
\end{aligned}
$$

which reduces to

$$
\begin{aligned}
& -a^{2} z^{n} \frac{\partial \omega}{\partial z^{n}} d z^{I}+\sum_{s}(-1)^{k+s} a^{2} z^{n} \frac{\partial \omega}{\partial \bar{z}^{i} s} d z^{I \backslash i} s{ }^{n} \Lambda z^{n} \\
& \text { if } i_{k}<n \\
& 0 \\
& \text { if } \quad i_{k}=n
\end{aligned}
$$

(4) $\rho L_{\nabla \rho} \omega=$

$$
\begin{gathered}
\left(a \bar{z}^{n}+b\right)\left(\nabla \rho\left(\omega_{I}\right) d \bar{z}^{I}+\omega_{I} \frac{\partial \rho^{i} s}{\partial \bar{z}^{j}} d \bar{z}^{i} 1^{\left.\cdots i_{s-1} j i_{s+1} \cdots{ }^{i_{k}}\right)}\right. \\
=a^{2} \bar{z}^{n} \frac{\partial \omega_{I}}{\partial \bar{z}^{-n}} d^{-I}+\cdots
\end{gathered}
$$

which remains as $a^{2} z^{n} \frac{\partial \omega}{\partial z^{n}} d z^{I}$ for any $I$


$$
=\left(\mathrm{a} \overline{\mathrm{z}}^{\mathrm{n}}+\mathrm{b}\right) \rho^{\mathrm{i}} \frac{\partial \omega_{\mathrm{I}}}{\partial \bar{z}^{j}} \iota \frac{}{\partial / \partial \bar{z}^{i}}\left(\mathrm{~d} \bar{z}^{j} \Lambda \mathrm{~d} \overline{\mathrm{z}}^{\mathrm{I}}\right)
$$

$$
=a^{2} z^{n} \frac{\partial \omega_{I}}{\partial \bar{z}^{j}} \iota_{\partial / \partial \bar{z}^{n}}\left(d^{-j} \Lambda d \bar{z}^{I}\right)+\cdots
$$

recalling $\rho^{i}=O(\rho)$ if $i<n$
yielding

$$
\begin{gathered}
a^{2} z^{n} \frac{\partial \omega}{\partial z^{n}} d z^{I} \quad \text { if } \quad i_{k}<n \\
a^{2} \sum(-1)^{k} z^{n} \frac{\partial \omega_{I}}{\partial z^{j}} d z^{j} \wedge d z^{I \backslash n} \quad \text { if } \quad i_{k}=n
\end{gathered}
$$

(6)

$$
\begin{array}{r}
\mathrm{d} \rho \wedge \iota_{\nabla \rho} \omega=\left(\operatorname{ad} \overline{\mathbf{z}}^{\mathrm{n}}+\cdots\right) \wedge \rho^{\mathrm{i}}{ }_{\left(\partial / \partial \bar{z}^{\mathrm{i}}\right)}{ }^{\omega} \mathrm{I} \mathrm{~d}^{-\mathrm{I}} \\
=\mathrm{a}^{2} \mathrm{~d}^{-\mathrm{z}} \wedge \omega_{I^{\iota}}^{\left(\partial / \partial \bar{z}^{-\mathrm{n}}\right)} \mathrm{d}^{-\mathrm{I}}
\end{array}
$$

and all that remains is

$$
\begin{gathered}
a^{2} \omega_{I} d z^{I} \quad i_{k}=n \\
0 \quad i_{k}<n
\end{gathered}
$$

(7) Finally, $\mathrm{d} \rho(\nabla \rho) \omega=\mathrm{a}^{2} \omega_{\mathrm{I}} \mathrm{d} \overline{\mathbf{z}}^{\mathrm{I}}+\cdots$
leaves $\quad a^{2} \omega_{I} d z^{I} \quad$ for any $\quad I$.

Heuristically, the procedure throughout is to expand the coefficients of $P$ in a power series, discard any term in which the power of $\bar{z}^{\text {n }}$ exceeds the number of derivatives, and then replace all coordinates $\bar{z}^{i}$ by their linear analogues $z^{i}$.

Collect all terms above with the appropriate coefficients to obtain

$$
\begin{aligned}
& \left\{a ( 0 ) ^ { 2 } \left\{\left[-\left(z^{n}\right)^{2} \sum_{i=1}^{n} \frac{\partial^{2} \omega}{\left(\partial z^{i}\right)^{2}}+\right.\right.\right. \\
& \text { (3.2) }
\end{aligned}
$$

Furthermore, the product formula (2.10) shows that

$$
N_{p}\left(\Delta_{g}\right)=\left(z^{n}\right)^{-k_{N}}(P)\left(z^{n}\right)^{k}
$$

so that $N_{p}\left(\Delta_{g}\right)\left[\frac{\omega}{\left(z^{n}\right)^{k}}\right]=\frac{N_{p}(P) \omega}{\left(z^{n}\right)^{k}}$, i.e. the action induced by $N_{p}\left(\Delta_{g}\right)$ on ${ }^{0} \Lambda^{k}\left(M_{p}\right)$ is given by $N_{p}(P)$, as claimed.

Finally, to see that $N_{p}\left(\Delta_{g}\right)$ is the Laplacian of the metric $\left(a(0) z^{n}\right)^{-2} d z^{2}$, simply observe that this a special case of all of the calculations above, and preceding (1.21), with $\rho=a(0) z^{n}, \quad h=d z^{2}$.

The corresponding $P$ is already dilation invariant and so no terms are lost in passing to the limit in (2.8). The assertion now follows, for (3.2) depends only on the dimension $n$, the degree $k$, and $\left|d \rho_{p}\right|_{n}^{2}$.

In the next section we construct an inverse for the action of $N_{p}\left(\Delta_{g}\right)$ on ${ }^{o} \Lambda^{k}\left(M_{p}\right)$, i.e. for $N_{p}(P)$. However, it is really the normal operator, as defined by (2.43), acting on sections of $\left.{ }^{o} \Lambda_{\ell}^{k} \otimes \Gamma_{0}^{1 / 2}\right|_{F_{p}}$ that we need to invert. It is true, though not immediately obvious, that this bundle and ${ }^{0} \Lambda^{k}\left(M_{p}\right) \otimes \Gamma_{0}^{1 / 2}\left(M_{p}\right)$ are naturally isomorphic. Thus it suffices to study the operator over $M_{p}$, then transfer the results back to $F_{p}$. Let us now examine the details of this bundle isomorphism. We prove the identification only for the $k$-form bundle; a similar argument applies equally well to the half-densities. In fact, the density bundles will be systematically neglected through most of this chapter. Their later inclusion is effected by 'conjugating the whole discussion by $x^{n / 2}$.
(3.3) Lemma: Both ${ }^{o} \Lambda^{k}\left(M_{p}\right)$ and $\left.{ }^{o} \Lambda_{\ell}^{k}\right|_{F_{p}}$ are trivial, though not naturally so. There is a canonical equivalence between these two bundles over the interiors of their bases. Furthermore, each is also naturally equivalent to ${ }^{o} \Lambda^{k}\left(B^{n}\right)=\Lambda^{k}\left({ }^{o} T^{*} B^{n}\right)$, again only over the interiors of their bases.

Proof: ${ }^{o} \Lambda^{k}\left(M_{p}\right)$ lifts to the left factor of $M_{p} \times M_{p}$, thence to $M_{p} \times_{0} M_{p}$; the resulting bundle we denote ${ }^{o} \Lambda_{\ell}^{k}\left(M_{p} \times{ }_{0} M_{p}\right)$. It then pulls back to the fibre $F_{0}$ over $(0,0) \in \partial M_{p} \times \partial M_{p}$. Examination of the rules (2.25) for formal pullback show that this restriction is isomorphic to ${ }^{o} \Lambda^{k}\left(M_{p}\right)$. Now, the linear model $M_{p}$ is equivalent to first order with $M$ at $p$. This means that the stretched products $M_{p} \times_{0} M_{p}$ and $M \times{ }_{0} M$ are also first order equivalent near $F_{0}$ and $F_{p}$, respectively. Hence ${ }^{o} \Lambda^{k}\left(M_{p}\right)\left(M_{p} \times_{0} M_{p}\right)$ and ${ }^{o} \Lambda_{\ell}^{k}\left(M_{0} M^{M}\right)$ agree to order zero at these fibres, so their restrictions to $F_{0}$ and $F_{p}$ coincide in a natural manner. Finally, Lemma (2.19) asserts the equivalence of any of these bundles with ${ }^{o} \Lambda^{k}\left(B^{n}\right)$, at least over the interiors of the bases. Notice that the isomorphism ${ }^{o} \Lambda^{k}\left(M_{p}\right) \simeq{ }^{o} \Lambda^{k}\left(B^{n}\right)$ is not induced by the usual conformal map $\mathbb{R}_{+}^{\mathrm{n}} \longrightarrow \mathrm{B}^{\mathrm{n}}$.

As we have indicated earlier, the indicial operator of $P$ is an ordinary differential operator of Euler type; it
is actually an uncoupled system. Using the coordinates of (3.1) with $p$ corresponding to 0 , from (3.2) we derive
(3.4) $\quad I_{p}(P)\left(\omega_{I} d z^{I}\right)=\{$

$$
\begin{aligned}
& a(0)^{2}\left[-\left(z^{n}\right)^{2} \frac{\partial^{2} \omega_{I}}{\left(\partial z^{n}\right)^{2}}+(n-2) \frac{\partial \omega_{I}}{\partial z^{n}}\right. \\
& \left.-(k-1)(n-k) \omega_{I}\right] d z^{I} \quad \text { if } \quad i_{k}=n
\end{aligned}
$$

Reverting now to the notation $z=(y, x)$, the solutions of

$$
I_{p}(P)(\omega)=0
$$

are of the form $\omega=\mathrm{x}^{\mathrm{s}}{ }^{1} \mathrm{dy}{ }^{\mathrm{I}}, \quad \omega=\mathrm{x}^{\mathrm{S}}{ }^{2} \mathrm{dy}{ }^{\mathrm{J}} \wedge \mathrm{dx}, \quad|\mathrm{I}|=\mathrm{k}$, $|J|=k-1$. The numbers $s_{1}$ and $s_{2}$ are called the indcal roots. A short calculation shows that

$$
s_{1}=k \quad \text { or } \quad n-k-1 \quad \text { unless } \quad k=\frac{n-1}{2}
$$

(3.5)

$$
s_{2}=k-1 \text { or } n-k \quad \text { unless } \quad k=\frac{n+1}{2} .
$$

The solutions in the exceptional cases $k=(n \neq 1) / 2$ involve additional logarithmic factors, but we are not concerned with these here.

The basic reason the parametrix construction breaks down in these exceptional cases is that only for these values of $k$ do the indicial roots not vanish at different rates at $x=0$. For the other values of $k$, only one of the solutions corresponding to each pair of indicial roots lies in the $L^{2}$ space which is natural for the problem. Thus:

Both $x^{n-k}, x^{n-k-1} \in L_{\text {loc }}^{2}\left(x^{-n} d x\right)$ iff $k<\frac{n-1}{2}$

$$
\begin{equation*}
\text { Both } x^{k}, x^{k-1} \in L_{l o c}^{2}\left(x^{-n} d x\right) \quad \text { iff } k>\frac{n+1}{2} \tag{3.6}
\end{equation*}
$$

Slightly restated, if $\omega$ solves $I_{p}(P) \omega=0$ and $\omega \epsilon$ $L_{\text {loc }}^{2}\left(x^{-n} d x\right)$, then (with summation intended)

$$
\omega=c_{I} x^{n-k-1} d y{ }^{I}+c_{J} x^{n-k} d y{ }^{J} \wedge d x \text { if } k<\frac{n-1}{2}
$$

$$
\begin{equation*}
\omega=c_{I} x^{k} d y{ }^{I}+c_{J} x^{k-1} d y y^{J} \wedge d x \quad \text { if } k>\frac{n+1}{2} \tag{3.7}
\end{equation*}
$$

We shall prove in the next chapter that if $\omega \in L^{2}(\mathrm{dg})$, $\mathrm{P} \omega=0$, then it is polyhomogeneous conormal at $\partial \mathrm{M}$, and (3.7) is the leading term in its asymptotic expansion there.

We should remark, as a general note, that the anaylsis of this particular $V_{0}$ elliptic problem is greatly simplified by the fact that the indicial roots (3.5) are independent of the base point $p$. In the general problem they may well vary, and complications arise from repeated roots, etc.
B. The Inverse of the Hyperbolic Laplacian

Let $N_{P}$ denote the operator of (3.2), neglecting the a(0) ${ }^{2}$ factor. From Lemma (3.3), $N_{P}$ may be interpreted as the operator induced by the constant curvature Laplacian either on ${ }^{o} \Lambda^{k}\left(\mathbb{R}_{+}^{n}\right)$ or ${ }^{o} \Lambda^{k}\left(B^{n}\right)$. It is important that we have both of these models available. The inverse for $N_{P}$ is constructed using the Fourier transform on the $y$ variables in the upper half space. This is quite natural, for these hypersurfaces are horospheres, hence have natural Euclidean structures. Then, however, the additional rotational symmetry of the ball is useful in understanding the behaviour of this inverse near $\infty \in \mathbb{R}_{+}^{n}$.

So, $N_{P}$ has a partial symbol given by

$$
\mathrm{N}_{\mathrm{P}}\left(\mathrm{x} \partial_{\mathrm{x}}, \mathrm{x} \partial_{\mathrm{y}}\right) \int \mathrm{e}^{\mathrm{iy} \mathrm{\eta} \hat{\omega}(\eta) \mathrm{d} \eta=\int \mathrm{e}^{i y \cdot \eta} \cdot \mathrm{~N}_{\mathrm{P}}\left(\mathrm{x} \partial_{\mathrm{x}}, \mathrm{ix} \eta\right) \hat{\omega}(\eta) \mathrm{d} \eta}
$$

where $\eta=\eta_{1} \mathrm{dy}^{1}+\cdots+\eta_{\mathrm{n}-1} \mathrm{dy}^{\mathrm{n}-1}$ is dual to y , and

$$
\hat{\omega}(\eta)=\int e^{-i y \cdot \eta} \omega(y) d y
$$

is the usual Fourier transform on the components of $\omega$. It is easy to see from (3.2) that

$$
\begin{aligned}
& N_{P}\left(x \partial_{x}, i x \eta\right) \hat{\omega}_{I} d y \\
&= {\left[-x^{2} \frac{\partial^{2} \hat{\omega}_{I}}{\partial x^{2}}+(n-2) x \frac{\partial \hat{\omega}_{I}}{\partial x}\right.} \\
&\left.+\left(x^{2}|\eta|^{2}-k(n-k-1)\right) \hat{\omega}_{I}\right] d y{ }^{I} \\
&+2(-1)^{k_{\iota}(i x \eta) \hat{\omega}_{I} d y}{ }^{I} \wedge d x
\end{aligned}
$$

(3.8)

$$
\begin{aligned}
N_{P}\left(x \partial_{x}, i x \eta\right) \hat{\omega}_{J} d y & \wedge d x=\left[-x^{2} \frac{\partial^{2} \hat{\omega}_{J}}{\partial x^{2}}+(n-2) x^{\frac{\partial \hat{\omega}}{\partial x}}\right. \\
& \left.+\left(x^{2}|\eta|^{2}-(k-1)(n-k)\right) \hat{\omega}_{J}\right] d y{ }^{J} \wedge d x \\
& +2(-1)^{k+1} i x \eta \wedge \hat{\omega}_{J} d y J
\end{aligned}
$$

The $\tilde{\eta}$ in the first of these formulae is the $h$-dual of $\eta$ :

$$
\tilde{\eta}=\eta_{1} \frac{\partial}{\partial y^{1}}+\cdots+\eta_{\mathrm{n}-1} \frac{\partial}{\partial \mathrm{y}^{\mathrm{n}-1}}
$$

and is to be contracted with $d y{ }^{I} \wedge d x$.
The operator of (3.8), which we also call $N_{P}$, is only mildly coupled. This allows all solutions to $N_{P} \hat{\omega}=0$
to be found, whence the inverse may be constructed and analyzed. Thus, regarding $\eta$ now as a fixed parameter, the $y$ coordinates may be rotated so that (3.8) assumes a particularly simple form. In fact, rotate $y$ so that $\eta=$ $|\eta| \mathrm{dy}^{1}$. Then

$$
\begin{gathered}
\eta \wedge \mathrm{dy}^{\mathrm{J}} \neq 0 \Leftrightarrow \mathrm{j}_{1}>1 \\
\iota(\tilde{\eta}) \mathrm{dy}{ }^{\mathrm{I}} \wedge \mathrm{dx} \neq 0 \Leftrightarrow \mathrm{i}_{1}=1
\end{gathered}
$$

There are three cases to consider:

$$
\begin{equation*}
\text { i) } \hat{\omega}=\hat{\omega}_{\mathrm{I}} \mathrm{dy}^{\mathrm{I}}, \quad \mathrm{i}_{1}>1 . \tag{3.9}
\end{equation*}
$$

$$
\Rightarrow N_{P} \hat{\omega}=\left[-x^{2} \frac{\partial^{2} \hat{\omega}_{I}}{\partial x^{2}}+(n-2) x \frac{\partial \hat{\omega}_{I}}{\partial x}+\left(x^{2}|\eta|^{2}-k(n-k-1)\right) \hat{\omega}_{I}\right] d y{ }^{I}
$$

ii) $\quad \hat{\omega}=\hat{\omega}_{\mathrm{J}} \mathrm{dy} \mathrm{J}^{\mathrm{J}} \mathrm{dx}, \quad \mathrm{j}_{1}=1$.

$$
\Rightarrow N_{P} \hat{\omega}=\left[-x^{2} \frac{\partial^{2 \hat{\omega}} J}{\partial x^{2}}+(n-2) x \frac{\partial \hat{\omega} J}{\partial x}\right.
$$

$$
\left.+\left(x^{2}|\eta|^{2}-(k-1)(n-k)\right) \hat{\omega}_{J}\right] d y^{J} \wedge d x
$$

iii) $\hat{\omega}=\hat{\omega}_{\mathrm{I}} \mathrm{dy}{ }^{\mathrm{I}}+\hat{\omega}_{\mathrm{J}} \mathrm{dy}{ }^{\mathrm{J}} \wedge \mathrm{dx}$,

$$
I=\left(i_{1}, \cdots, i_{k}\right)=\left(1, j_{1}, \cdots, j_{k-1}\right)=(1, J)
$$

$$
\begin{aligned}
\Rightarrow N_{P} \hat{\omega}=\left[-x^{2} \frac{\partial^{2} \hat{\omega}_{I}}{\partial x^{2}}+(n-2)\right. & x \frac{\partial \hat{\omega}_{I}}{\partial x}+\left(x^{2}|\eta|^{2}-k(n-k-1)\right) \hat{\omega}_{I} \\
& \left.+2(-1)^{k+1} i x|\eta| \hat{\omega}_{J}\right] d y y^{I} \\
+ & {\left[-x^{2} \frac{\partial^{2} \hat{\omega} J}{\partial x^{2}}+(n-2) x \frac{\partial \hat{\omega}_{J}}{\partial x}+\left(x^{2}|\eta|^{2}-(k-1)(n-k)\right) \hat{\omega}_{J}\right.} \\
& +2(-1)^{\left.k_{i x}|\eta| \hat{\omega}_{I}\right] d y}{ }^{J} \wedge d x
\end{aligned}
$$

The first two of these operators are quite easy to invert. The third one is slightly less so, since not all of its solutions are known explicitly. Nonetheless, it may be regarded as a compact perturbation of the diagonal system formed by i) and ii).

The program now is as follows. Ultimately we seek to prove that the partial differential operator $N_{P}$ of (3.2), acting on

$$
L^{2}\left(\mathbb{R}^{+} \times \mathbb{R}^{n-1}, x^{-n} d x d y\right) \simeq L^{2}\left(\mathbb{R}^{+}, x^{-n} d x ; L^{2}\left(\mathbb{R}^{n-1}, d y\right)\right)
$$

has a bounded inverse. Yet more refined regularity properties of this inverse are proved in the next section. Upon conjugation by the Fourier transform it suffices to show that the ordinary differential operator $N_{P}$ of (3.8), acting on

$$
L^{2}\left(\mathbb{R}^{+}, x^{-n} d x ; L^{2}\left(\mathbb{R}^{n-1}, d \eta\right)\right)
$$

also has a bounded inverse.
(3.10) Proposition: There exists a unique kernel $G(x, \tilde{x}, \eta)$ which is bounded on

$$
L^{2}\left(\mathbb{R}^{+}, \mathrm{x}^{-\mathrm{n}} \mathrm{dx} ; \mathrm{L}^{2}\left(\mathbb{R}^{\mathrm{n}-1}, \mathrm{~d} \eta\right)\right)
$$

and such that

$$
N_{P}\left(x \partial_{x}, i x \eta\right) G(x, \tilde{x}, \eta)=x^{n} \delta(x-\tilde{x})
$$

(The right side of this equation is the Schwartz kernel of the identity on the space above.)

The proof is quite long and involves several lemmas along the way. Consider first the two operators (3.9) i), ii). Set

$$
\begin{aligned}
& P_{1}=-x^{2} \frac{d^{2}}{d x^{2}}+(n-2) x \frac{d}{d x}+\left(x^{2}|\eta|^{2}-k(n-k-1)\right) \\
& P_{2}=-x^{2} \frac{d^{2}}{d x^{2}}+(n-2) x \frac{d}{d x}+\left(x^{2}|\eta|^{2}-(k-1)(n-k)\right)
\end{aligned}
$$

These are Fuchsian operators, the indicial roots of which, as expected, are $k, n-k-1$ and $k-1, n-k$. In fact, for $k<\frac{n-1}{2}$ (the case $k>\frac{n+1}{2}$ is treated similarly, and
shall be mentioned only rarely) we find, cf. [19],

$$
\begin{gathered}
\mathrm{P}_{1} \mathrm{u}_{1}=0 \Rightarrow \mathrm{u}_{1}(\mathrm{x})=\mathrm{c}_{1} \mathrm{x}^{\frac{\mathrm{n}-1}{2}} \mathrm{I}_{v_{1}}(\mathrm{x}|\eta|)+\mathrm{c}_{2} \mathrm{x}^{\frac{\mathrm{n}-1}{2}} \mathrm{~K}_{v_{1}}(\mathrm{x}|\eta|) \\
v_{1}=\frac{\mathrm{n}-2 \mathrm{k}-1}{2}
\end{gathered}
$$

(3.11)

$$
\begin{gathered}
P_{2} u_{2}=0 \Rightarrow u_{2}(x)=c_{1} x^{\frac{n-1}{2}} I_{v_{2}}(x|\eta|)+c_{2} x^{\frac{n-1}{2}} K_{v_{2}}(x|\eta|) \\
v_{2}=\frac{n-2 k+1}{2} .
\end{gathered}
$$

The $I_{\nu}$ amd $K_{v}$ are Bessel functions -specifically, the modified Bessel function of the first kind and Macdonald's function, respectively, of order $v>0$. Their asymptotics are well-known [19]

$$
\begin{array}{r}
\mathrm{I}_{v}(\mathrm{x}) \sim \mathrm{x}^{v} / 2^{v} \Gamma(1+v), \quad \mathrm{K}_{v}(\mathrm{x}) \sim 2^{v-1} \Gamma(v) / \mathrm{x}^{v} \\
\mathrm{x} \longrightarrow 0^{+}
\end{array}
$$

(3.12)

$$
\begin{array}{r}
I_{v}(x) \sim e^{x} / \sqrt{2 \pi x}, \quad K_{v}(x) \sim e^{-x \sqrt{\pi / 2 x}} \\
x \longrightarrow \infty
\end{array}
$$

Thus, near $\mathbf{x}=0$,

$$
\mathrm{x}^{\frac{\mathrm{n}-1}{2}} \mathrm{I}_{v_{1}}(\mathrm{x}|\eta|) \sim \mathrm{x}^{\mathrm{n}-\mathrm{k}-1}|\eta|^{v_{1}}
$$

$$
\mathrm{x}^{\frac{\mathrm{n}-1}{2}} \mathrm{~K}_{v_{1}}(\mathrm{x}|\eta|) \sim \mathrm{x}^{\mathrm{k}}|\eta|^{-v_{1}}
$$

(3.13)

$$
\begin{aligned}
& x^{\frac{\mathrm{n}-1}{2}} \mathrm{I}_{\nu_{2}}(\mathrm{x}|\eta|) \sim \mathrm{x}^{\mathrm{n}-\mathrm{k}}|\eta|^{v_{2}} \\
& \\
& \mathrm{x}^{\frac{\mathrm{n}-1}{2} \mathrm{~K}_{\nu_{2}}(\mathrm{x}|\eta|) \sim \mathrm{x}^{\mathrm{k}-1}|\eta|^{-v_{2}}}
\end{aligned}
$$

(up to constant factors). Recalling (3.7), $u_{i} \epsilon$ $L_{\text {loc }}^{2}\left(x^{-n} d x\right)$ only if $c_{2}=0, i=1,2$. On the other hand, these solutions are exponentially increasing at infinity. The other pair, with $c_{1}=0$, decrease suitably at infinity, but do not vanish to sufficiently high order at $x=0$. (For the exceptional degrees $k=(n+1) / 2$, the $I_{v}$ and $K_{v}$ differ only by a logarithmic factor at $x=$ 0.)

Now we construct kernels inverting $P_{1}$ and $P_{2}$. In fact, these two operators are treated so similarly that only the work for $\mathrm{P}_{1}$ is displayed. With

$$
u(x)=x^{\frac{n-1}{2}} I_{v_{1}}(x|\eta|), \quad v(x)=x^{\frac{n-1}{2}} K_{v_{1}}(|\eta|)
$$

we know the kernel is of the form

$$
G_{1}(x, \tilde{x}, \eta)=\alpha(\tilde{x}) u(x) H(\tilde{x}-x)+\beta(\tilde{x}) v(x) H(x-\tilde{x})
$$

and satisfies

$$
P_{1} G_{1}(x, \tilde{x}, \eta)=x^{n} \delta(x-\tilde{x})
$$

which is the kernel of the identity on $L^{2}\left(x^{-n} d x\right)$. This form is chosen so that $G_{1}$ is a 'free' solution for $x \neq \tilde{x}$, satisfying the boundary conditions (with respect to x) of sufficiently rapid decay at $x=0, \infty$. A bit of calculation, with $\alpha(\tilde{x})=\gamma(\tilde{x}) v(\tilde{x}), \beta(\tilde{x})=\gamma(\tilde{x}) u(\tilde{x})$, shows

$$
P_{1} G_{1}=x^{n} \delta(x-\tilde{x})\left[-2 \gamma(x) x^{2-n}\left(u(x) v^{\prime}(x)-u^{\prime}(x) v(x)\right]\right.
$$

Now, as $P_{1}=-x^{n} \frac{d}{d x}\left(x^{2-n} \frac{d}{d x}\right)+\left(x^{2}|\eta|^{2}-k(n-k-1)\right)$, the expression

$$
x^{2-n}\left(u(x) v^{\prime}(x)-u^{\prime}(x) v(x)\right)
$$

is simply the Sturm-Liouville Wronskian of $\mathrm{P}_{1}$, hence independent of $x$. By examining, for example, its power series as $x \longrightarrow 0, ~ i t ~ i s ~ a l s o ~ i n d e p e n d e n t ~ o f ~ \eta . ~ C h o o s-~$ ing $r(x)$ to be the appropriate constant
(3.14)

$$
\begin{aligned}
G_{1}(x, \tilde{x}, \eta)= & c_{k, n} x^{\frac{n-1}{2} \tilde{x}^{\frac{n-1}{2}}}\left[I_{v_{1}}(x|\eta|) K_{v_{1}}(\tilde{x}|\eta|) H(\tilde{x}-x)\right. \\
& \left.+I_{v_{1}}(\tilde{x}|\eta|) K_{v_{1}}(x|\eta|) H(x-\tilde{x})\right]
\end{aligned}
$$

Analogously, the kernel $G_{2}$ for $P_{2}$ is obtained by replacing $v_{1}$ with $v_{2}$.

It is convenient to introduce the space
$(3.15) \quad H_{r}^{(p)}=\hat{H}^{p}\left([0,2] ; x^{-n} d x\right)+x^{r} H^{p}([1, \infty) ; d x)$
where $p \in \mathbb{Z}^{+}, r \in \mathbb{R}, H^{p}$ is the usual Sobolev space of order $p$ and

$$
\hat{H}^{p}=\left\{f: x^{i} \frac{d^{j} f}{d x^{j}} \in L^{2}\left([0,2] ; x^{-n} d x\right) \text { for } p \geq i \geq j \geq 0\right\}
$$

so as to formulate the
(3.16) Proposition: For $k<\frac{n-1}{2}$ or $k>\frac{n+1}{2}$

$$
P_{i}: \#_{\mathrm{r}}^{(2)} \longrightarrow H_{\mathrm{r}+2}^{(0)} \quad i=1,2
$$

is an isomorphism, with inverse represented by the kernel $G_{i}(x, \tilde{x}, \eta)$.

Proof: Obviously $P_{1}$ and $P_{2}$ map $\mathscr{H}_{\mathrm{r}}^{(2)}$ into $\mathscr{H}_{\mathrm{r}+2}^{(0)}$, so it suffices to show that

$$
\mathrm{G}_{\mathrm{i}}: \mathscr{H}_{\mathrm{r}+2}^{(0)} \longrightarrow \mathscr{H}_{\mathrm{r}}^{(2)} \quad \mathrm{i}=1,2
$$

is bounded. By construction it must then actually invert $P_{i}$. Uniqueness is also immediate since $P_{i} u=0$ has no solutions in $\mathscr{H}_{\mathrm{r}}^{(2)}$ for any $r$.

The basic tool to prove boundedness of a kernel is

## Schur's Test:

The positive kernel $G(x, \tilde{x})$ induces a bounded map on $L^{2}(d \mu)$ if there exist positive measurable functions $p$ and $q$ such that

$$
\int \mathrm{G}(\mathrm{x}, \tilde{\mathrm{x}}) \mathrm{q}(\tilde{\mathrm{x}}) \mathrm{d} \mu(\tilde{\mathrm{x}}) \leq \alpha \mathrm{p}(\mathrm{x}) ; \quad \int \mathrm{G}(\mathrm{x}, \tilde{\mathrm{x}}) \mathrm{p}(\mathrm{x}) \mathrm{d} \mu(\mathrm{x}) \leq \mathrm{Bq}(\tilde{\mathrm{x}})
$$

Now, given the decomposition of $\mathscr{H}_{\mathrm{r}+2}^{(0)}$ as a sum in (3.14) it suffices to prove boundedness on each summand. Thus if $\mathrm{f} \in \mathscr{H}_{\mathrm{r}+2}^{(0)}$, write $\mathrm{f}=\mathrm{f}_{1}+\mathrm{f}_{2}, \mathrm{f}_{1} \in \mathrm{~L}^{2}\left([0,1] ; \mathrm{x}^{-\mathrm{n}} \mathrm{dx}\right)$, $f_{2} \in x^{r+2} L^{2}([1, \infty) ; d x)$. Accordingly

$$
\|f\|_{r,(0)}^{2}=\int_{0}^{1}\left|f_{1}(x)\right|^{2} x^{-n} d x+\int_{1}^{\infty}\left(x^{-r-2} f_{2}(x)\right)^{2} d x
$$

Case I: $f \in L^{2}\left([0,1], x^{-n} d x\right)=x^{n / 2} L^{2}([0,1], d x)$.

It is convenient to include the various weight factors in the kernel. Thus let

$$
J(x, \tilde{x})=x^{-n / 2} G_{1}(x, \tilde{x}) \tilde{x}^{-n / 2}
$$

Then $G_{1}: L^{2}\left([0,1], x^{-n} d x\right) \longrightarrow \hat{H}^{2}\left([0,1], x^{-n} d x\right)$ is bounded if and only if

$$
\begin{aligned}
& L^{2}(d x) \ni f \longmapsto x^{n / 2} f \longmapsto \int G_{1}(x, \tilde{x}) \tilde{x}^{n / 2} f(\tilde{x}) \tilde{x}^{-n} d \tilde{x} \\
& \longmapsto \int x^{-n / 2} G_{1}(x, \tilde{x}) \tilde{x}^{-n / 2} f(\tilde{x}) d \tilde{x} \\
&=\int J(x, \tilde{x}) f(\tilde{x}) d \tilde{x} \in \hat{H}^{2}([0,1] ; d x)
\end{aligned}
$$

is also bounded. Furthermore, it even suffices to prove boundedness of $J$ on $L^{2}(d x)$ since $x d J / d x, x^{2} d^{2} J / d x{ }^{2}$ are similar to $J$ in their asymptotics as $x \longrightarrow 0$ and $x \longrightarrow \infty$. (These derivatives contain terms with $\delta(x-\tilde{x})$ as a factor, but never $\delta^{\prime}(x-\tilde{x})$, multiplied by a bounded coefficient; these are obviously bounded on $L^{2}$.) Now, since $x$ and $\tilde{x}$ lie in $[0,1]$ we may replace J by

$$
\begin{aligned}
J_{a}(x, \tilde{x}) & =x^{-n / 2}\left[x^{n-k-1} \tilde{x}^{2} H(\tilde{x}-x)+x^{k_{x}^{\sim}}{ }^{n-k-1} H(x-\tilde{x})\right] \tilde{x}^{-n / 2} \\
& =x^{\frac{n}{2}-k-1 \sim k-\frac{n}{2}} H(\tilde{x}-x)+x^{k-\frac{n}{2} \sim} x^{\frac{n}{2}-k-1} H(x-\tilde{x}) .
\end{aligned}
$$

For $k<\frac{n}{2}-1$, apply Schur's test with $p=q=1$ :

$$
\begin{aligned}
& \int_{0}^{1} J_{a}(x, \tilde{x}) d \tilde{x}=\int_{0}^{x} x^{k-\frac{n}{2} \tilde{x}^{\frac{n}{2}-k-1}} d \tilde{x}+\int_{x}^{1} x^{\frac{n}{2}-k-1 \sim \tilde{x}^{k-\frac{n}{2}}} d \tilde{x} \\
& =\left.x^{k-\frac{n}{2}} \cdot \frac{\tilde{x}^{\frac{n}{2}}-k}{\left(\frac{n}{2}-k\right)}\right|_{0} ^{x}+\left.x^{\frac{n}{2}-k-1} \frac{\tilde{x}^{k-\frac{n}{2}}}{\left(k-\frac{n}{2}+1\right)}\right|_{x} ^{1} \\
& =x^{k-\frac{n}{2}}\left(x^{\frac{n}{2}-k} /\left(\frac{n}{2}-k\right)\right)+x^{\frac{n}{2}-k-1}\left(1-x^{k-\frac{n}{2}+1}\right) /\left(k-\frac{n}{2}+1\right) \leq C
\end{aligned}
$$

since neither $k-\frac{n}{2}$ nor $\frac{n}{2}-k-1$ equal -1 , and $\frac{\mathrm{n}}{2}-\mathrm{k}-1>0$.

$$
\text { Symmetrically } \int_{0}^{1} J(x, \tilde{x}) d x \leq C
$$

For $k=\frac{n}{2}-1$ (the only excluded case), use $p=q=$ $(-\log x) / \sqrt{x}$.

$$
\begin{aligned}
& \int_{0}^{1} J_{a}(x, \tilde{x})(-\log \tilde{x}) / \tilde{x}^{1 / 2} d \tilde{x} \\
= & \int_{0}^{x} x^{-1} \frac{-\log \tilde{x}}{\tilde{x}} \tilde{x}^{1 / 2} \tilde{x}+\int_{x}^{1} \frac{-\log \tilde{x^{3 / 2}} \tilde{x}}{\tilde{x}^{3}} \tilde{x}
\end{aligned}
$$

$$
\begin{aligned}
& =\left.\mathrm{x}^{-1}\left(-2 \tilde{x}^{1 / 2} \log \tilde{x}+4 \tilde{x}^{1 / 2}\right)\right|_{0} ^{x}+\left.\left(2 \tilde{x}^{-1 / 2} \log \tilde{x}+4 \tilde{x}^{-1 / 2}\right)\right|_{x} ^{1} \\
& \leq C(-\log x) / \sqrt{x} \quad \text { and likewise for } \\
& \int_{0}^{1} J_{a}(x, \tilde{x})(-\log x) / \sqrt{x} d x .
\end{aligned}
$$

On the other hand, as $x \longmapsto \infty$ the only nonvanishing term in

$$
\int G_{1}(x, \tilde{x}) f(\tilde{x}) \tilde{x}^{-n} d \tilde{x}
$$

is, by (3.14)

$$
\int_{0}^{x} v(x) u(\tilde{x}) f(\tilde{x}) \tilde{x}^{-n} d \tilde{x}=\int_{0}^{1} v(x) u(\tilde{x}) f(\tilde{x}) \tilde{x}^{-n} d \tilde{x}
$$

since supp f $C[0,1]$. This decreases exponentially provided the integral makes sense. But $u \sim x^{n-k-1}$, so that

$$
\left.\begin{array}{l}
\int_{0}^{1} \tilde{x}^{n-k-1} f(\tilde{x}) \tilde{x}^{-n} d \tilde{x} \\
\quad \leq\left[\int_{0}^{1} f^{2} \tilde{x}^{\sim}-n\right. \\
d \tilde{x}
\end{array}\right]^{1 / 2}\left[\int_{0}^{1} \tilde{x}^{n-2 k-2} d \tilde{x}\right]^{1 / 2}<\infty . .
$$

Case II: $f \in x^{r+2} L^{2}([1, \infty) ; d x)$.

Again we combine weight factors into the kernel. Thus we need to prove boundedness of the composition
$\left.L^{2}[1, \infty) ; d x\right) \ni f \longmapsto x^{r+2} f \longmapsto \int G_{1}(x, \tilde{x}) \tilde{x}^{r+2-n_{f}}(\tilde{x}) d \tilde{x}$

$$
\longmapsto \int x^{-r} G_{1}(x, \tilde{x}) \tilde{x}^{r+2-n_{f}}(\tilde{x}) d \tilde{x}
$$

$$
\equiv \int J(x, \tilde{x}) f(\tilde{x}) d \tilde{x} \in H^{2}([1, \infty], d x)
$$

where $J(x, \tilde{x})=x^{-r} G_{1}(x, \tilde{x}) \tilde{x}^{r+2-n}$. As before it suffices to prove that $J$ is bounded on $L^{2}$, for its derivatives are handled similarly. We require the
(3.17) Lemma: If $\mu \in \mathbb{R}$ then as $x \longrightarrow \infty$,

$$
\begin{aligned}
& {\left[\int_{1}^{x} \tilde{x}^{\mu} \mathrm{e}^{\tilde{x}} \mathrm{dx}\right] \mathrm{x}^{-\mu} \mathrm{e}^{-\mathrm{x}} \longrightarrow 1} \\
& {\left[\int_{\mathrm{x}}^{\infty} \tilde{x}^{\mu} \mathrm{e}^{-\tilde{x}} \mathrm{~d} \tilde{\mathrm{x}}\right] \mathrm{x}^{-\mu} \mathrm{e}^{\mathrm{x}} \longrightarrow 1}
\end{aligned}
$$

Proof: Apply L'Hopital's rule.

Now by (3.12) we may replace $J$ by

$$
J_{a}(x, \tilde{x})=x^{\frac{n}{2}-r-1} \tilde{x}^{r+1-\frac{n}{2}}\left[e^{x-\tilde{x}_{H}(\tilde{x}-x)}+e^{\left.\tilde{x}-x_{H}(x-\tilde{x})\right]}\right.
$$

Apply Schur's test with $p=q=1$ :

$$
\begin{aligned}
& \int_{1}^{\infty} J_{a}(x, \tilde{x}) d \tilde{x}=\int_{1}^{x} x^{\frac{n}{2}-r-1} e^{-x} \tilde{x}^{r+1-\frac{n}{2}} e^{\tilde{x}} d \tilde{x} \\
& \quad+\int_{x}^{\infty} x^{\frac{n}{2}-r-1} e^{x} \tilde{x}^{r+1-\frac{n}{2}} e^{-\tilde{x}} d \tilde{x} \leq C \quad \text { by }(3.17) . \\
& \text { By symmetry, } \int_{1}^{\infty} J_{a}(x, \tilde{x}) d x \leq C .
\end{aligned}
$$

Finally, if $\operatorname{supp} f \subset[1, \infty]$, then as $x \longrightarrow 0$

$$
\begin{aligned}
& \int G_{1}(x, \tilde{x}) f(\tilde{x}) \tilde{x}^{-n} d \tilde{x}=\int_{x}^{\infty} u(x) v(\tilde{x}) f(\tilde{x}) \tilde{x}^{-n} d \tilde{x} \\
& =u(x) \int_{1}^{\infty} v(\tilde{x}) f(\tilde{x}) \tilde{x}^{-n} d \tilde{x} \in \hat{H}^{2}\left([0,1] ; x^{-n} d x\right)
\end{aligned}
$$

since the integral exists and is independent of $x$, and $u(x)$ lies in this space.

This concludes the proof of Proposition (3.16).

We must now construct an inverse for the operator of (3.9) iii). Set

$$
L_{0}=\left[\begin{array}{cc}
P_{1} & 0 \\
0 & P_{2}
\end{array}\right], \quad A=\left[\begin{array}{cc}
0 & 2(-1)^{k+1} i x|\eta| \\
2(-1)^{k_{i x}|\eta|} & 0
\end{array}\right]
$$

Then that operator is $L_{1}=L_{0}+A$. More generally, let

$$
L_{t}=L_{0}+t A, \quad 0 \leq t \leq 1
$$

By (3.16), for any r

$$
\mathrm{L}_{0}: \mathscr{H}_{\mathrm{r}}^{(2)} \longrightarrow \mathscr{H}_{\mathrm{r}+2}^{(0)}
$$

is invertible. (Strictly speaking, these spaces are defined to contain only scalar functions. Here of course we extend the definition to allow vector functions by requiring each component to lie in the space.) Clearly each $L_{t}$ maps $H_{r}^{(2)}$ into $H_{r+2}^{(0)}$.
(3.18) Lemma: $A: \mathscr{H}_{\mathrm{r}}^{(2)} \longrightarrow \mathscr{H}_{\mathrm{r}+2}^{(0)}$ is compact.

Proof: Certainly $A: x^{r} H^{2}([1, \infty), d x) \longrightarrow x^{r+1} H^{2}([1, \infty), d x)$ is bounded. The inclusion

$$
x^{r+1} H^{2} \longrightarrow X^{r+2} L^{2}
$$

is compact by the $L^{2}$ version of the Arzela-Ascoli Theorem. On the other hand, between

$$
\hat{\mathrm{H}}^{2}\left([0,2], \mathrm{x}^{-\mathrm{n}} \mathrm{dx}\right) \longrightarrow \mathrm{L}^{2}\left([0,2], \mathrm{x}^{-\mathrm{n}} \mathrm{dx}\right)
$$

A is the strong limit of the same multiplication operator acting between

$$
\hat{H}^{2}\left([\epsilon, 2], x^{-n} d x\right) \simeq \eta^{2}\left([\epsilon, 2], x^{-n} d x\right) \longrightarrow L^{2}\left([\epsilon, 2], x^{-n} d x\right)
$$

which, for each $\epsilon>0$, is compact as before. Such strong limits are compact provided a 'uniform smallness' condition at the boundary holds, cf. [29]. Here this states that if $f$ varies over a bounded set in $\hat{H}^{2}\left([0,2], x^{-n} d x\right)$, the functions Af have uniformly small norm in $L^{2}\left([0, \epsilon], x^{-n} d x\right)$. This follows from the factor of $x$ in A. The proof is complete.

By the Lemma,

$$
\mathrm{L}_{\mathrm{t}}: \mathscr{H}_{\mathrm{r}}^{(2)} \longrightarrow \mathscr{H}_{\mathrm{r}+2}^{(0)}
$$

is an analytic family of Fredholm operators. Thus, for any t :

$$
\begin{aligned}
\text { index } L_{t} & : \mathscr{H}_{\mathrm{r}}^{(2)} \longrightarrow \mathscr{H}_{\mathrm{r}+2}^{(0)} \\
& =\text { index } \mathrm{L}_{0}: \mathscr{H}_{\mathrm{r}}^{(2)} \longrightarrow \mathscr{H}_{\mathrm{r}+2}^{(0)}=0 .
\end{aligned}
$$

To prove that $L_{1}$ is an isomorphism we need only show
(3.20) Proposition: $L_{1} u=0, u \in \mathscr{H}_{r}^{(2)}$ for some $r \Rightarrow$ $\mathrm{u}=0$.
(3.21) Corollary: $L_{1}: \mathscr{H}_{\mathrm{r}}^{(2)} \longrightarrow \mathscr{H}_{\mathrm{r}+2}^{(0)}$ is invertible for each r.

Note that from (3.16), a solution to $L_{1} u=0$ lying in some $\mathscr{H}_{\mathbf{r}}^{(2)}$ must actually lie in every $\mathscr{H}_{\mathrm{r}}^{(2)},-\infty<r<\infty$. For
$u \in \mathscr{H}_{\mathrm{r}}^{(2)}, \quad\left(\mathrm{L}_{\mathrm{O}}+\mathrm{A}\right) \mathrm{u}=0 \Rightarrow \mathrm{~L}_{0} \mathrm{u}=-\mathrm{Au} \in \mathscr{H}_{\mathrm{r}+1}^{(0)} \Rightarrow \mathrm{u} \in \mathscr{H}_{\mathrm{r}-1}^{(2)}$
and the assertion follows inductively.

In essence, (3.20) is nothing but the statement that $\mathbb{H}^{n}$ has no $L^{2}$ harmonic forms outside the middle degrees. However, rather than developing this connection directly, which requires a fair bit of effort, it seems better to give a proof more nearly consonant with our overall techniques. There is actually an easy proof that $L_{1}$ is
strictly positive based on the inequality

$$
-x^{2} \frac{d^{2}}{d x^{2}}+(n-2) x \frac{d}{d x} \geq \frac{(n-1)^{2}}{4}
$$

-which is equivalent to the well-known fact that the spectrum of the hyperbolic Laplacian on functions is bounded below by $(n-1)^{2} / 4$. Unfortunately this fails in the two degrees $k=\frac{n}{2} \pm 1$.

The method which covers all cases involves an integration by parts to replace the second order system $N_{p} \hat{\omega}=0$ by the larger first order system-which may be solved-

$$
\mathrm{N}_{\mathrm{d}} \hat{\omega}=0, \quad \mathrm{~N}_{\delta} \hat{\omega}=0 .
$$

Obviously this is just the normal operator analogue of the fact that $L^{2}$ harmonic forms on $\mathbb{H}^{n}$ are both closed and coclosed.

Recall first that the conjugated hyperbolic Laplacian may be factored

$$
\begin{gathered}
N_{P}=x^{k} \Delta_{H} x^{-k} \\
=\left(x^{k} d x^{-k+1}\right)\left(x^{k-1} \delta_{H} x^{-k}\right)+\left(x^{k} \delta_{H} x^{-k-1}\right)\left(x^{k+1} d x^{-k}\right)
\end{gathered}
$$

Here $\delta_{H}$ is the coboundary operator on $\mathbb{H}^{n}$. Furthermore, by the product formula (2.10),

$$
N_{P}=N\left(x^{k} d x^{-k+1}\right) N\left(x^{k-1} \delta_{H} x^{-k}\right)
$$

$$
+N\left(x^{k} \delta_{H} x^{-k-1}\right) N\left(x^{k+1} d x^{-k}\right)
$$

(neglecting basepoints). Now, from (1.19) and (1.20) a calculation similar to that in Proposition (3.1) shows that the (conjugate by the Fourier transform of) the normal operators are

$$
N_{d} \hat{\omega}=(-1)^{k}\left(x \frac{\partial \hat{\omega}_{I}}{\partial x}-\hat{k}_{I}\right) d y{ }^{I} \wedge d x+i x \eta \wedge\left(\hat{\omega}_{I} d y{ }^{I}+\hat{\omega}_{J} d y{ }^{J} \Lambda d x\right)
$$

$$
\begin{equation*}
N_{\delta} \hat{\omega}=(-1)^{k}\left(x \frac{\partial \hat{\omega}_{J}}{\partial x}-(n-k) \hat{\omega}_{J}\right) d y J-i x \iota\left(\tilde{\eta}^{J}\right)\left(\hat{\omega}_{I} d y^{I}+\hat{\omega}_{J} d y{ }^{J} \Lambda d x\right) \tag{3.22}
\end{equation*}
$$

Using again the reduction $|\eta| d^{1}=\eta$, this system splits into three cases, just as in (3.9). Only the last of these possibilities is of interest here, and arises when $I=$ $(1, J)$ : then, for $\hat{\omega}=\hat{\omega}_{I} d y+\hat{\omega}_{J} d y{ }^{J} \wedge d x$

$$
\begin{equation*}
L_{d} \hat{\omega}=\left[(-1)^{k}\left(x \frac{d \hat{\omega}_{I}}{d x}-k \hat{\omega}_{I}\right)+i x|\eta| \hat{\omega}_{J}\right] d y{ }^{I} \wedge d x \tag{3.23}
\end{equation*}
$$

$$
L_{\delta} \hat{\omega}=\left[(-1)^{k}\left(x \frac{d \hat{\omega}}{d x}-(n-k) \hat{\omega}_{J}\right)-i x|\eta| \hat{\omega}_{I}\right] d y^{I}
$$

and the relationship

$$
L_{1}=L_{d} L_{\delta}+L_{\delta} L_{d}
$$

is still valid. (The degree $k$ in $L_{d}$ and $L_{\delta}$ has been suppressed, so in this last equation it is different in the two $L_{d}$ 's, and also in the two $L_{\delta}$ 's.)
(3.24) Lemma: Let $\hat{\omega}=f_{I} d y{ }^{I}+g_{J} d^{J} \Lambda d x$ and $\hat{\mu}=$ $\alpha_{K} d y^{K}+\beta_{L} d y^{L} \Lambda d x$ be $k$-form and ( $k-1$ )-form valued fundlions, respectively, on $\mathbb{R}^{+}$. Then using the pointwise Euclidean scalar product on forms we have

$$
\int_{0}^{\infty}\left\langle\mathrm{N}_{\mathrm{d}} \hat{\omega}, \hat{\mu}\right\rangle \mathrm{x}^{-\mathrm{n}} \mathrm{dx}=\int_{0}^{\infty}\left\langle\hat{\omega}, \mathrm{N}_{\delta} \hat{\mu}\right\rangle \mathrm{x}^{-\mathrm{n}} \mathrm{dx}
$$

if and only if for each $I: \quad \int_{0}^{\infty} \frac{d}{d x}\left(x^{1-n_{f}} \bar{\beta}_{I}\right) d x=0$.

Proof: $\int_{0}^{\infty}\left\langle N_{d} \hat{\omega}, \hat{\mu}\right\rangle x^{-n} d x$

$$
\begin{aligned}
& =\int_{0}^{\infty}\left\langle(-1)^{k}\left(x \frac{d f}{d x}-k f_{I}\right) d y^{I} \wedge d x+i x \eta \wedge \hat{\omega}, \hat{\mu}\right\rangle x^{-n} d x \\
& =\int_{0}^{\infty}\left[(-1)^{k}\left(x \frac{d f_{I}}{d x}-k f_{I}\right) \bar{\beta}_{I}+\langle\hat{\omega},-i x \iota(\tilde{\eta}) \hat{\mu}\rangle\right] x^{-n} d x
\end{aligned}
$$

$$
\begin{aligned}
& =\int_{0}^{\infty}(-1)^{\mathrm{k}}\left(-\mathrm{f}_{\mathrm{I}} \frac{\mathrm{~d}}{\mathrm{dx}}\left(\mathrm{x}^{\left.\left.1-\mathrm{n}_{\bar{\beta}_{\mathrm{I}}}\right)-\mathrm{k} \bar{\beta}_{\mathrm{I}} \mathrm{x}^{-\mathrm{n}}\right)-\langle\hat{\omega}, \mathrm{ix} \iota(\tilde{\eta}) \hat{\mu}\rangle \mathrm{x}^{-\mathrm{n}} \mathrm{dx}}\right.\right. \\
& +\int_{0}^{\infty}(-1)^{\mathrm{k}} \frac{\mathrm{~d}}{\mathrm{dx}}\left(\mathrm{x}^{1-\mathrm{n}_{\mathrm{f}}} \mathrm{I}_{\mathrm{I}}\right) \mathrm{dx} \\
& =\int_{0}^{\infty}\left\langle\hat{\omega}, \mathrm{N}_{\delta} \hat{\mu}\right\rangle \mathrm{x}^{-\mathrm{n}} \mathrm{dx}+\int_{0}^{\infty}(-1)^{\mathrm{k}} \frac{\mathrm{~d}}{\mathrm{dx}}\left(\mathrm{x}^{1-\mathrm{n}_{\mathrm{f}}} \mathrm{I}_{\mathrm{I}} \overline{\mathrm{~B}}_{\mathrm{I}}\right) \mathrm{dx} .
\end{aligned}
$$

Hence $N_{d}$ and $N_{\delta}$ are adjoint iffy this boundary term vanishes.

Proof of Proposition (3.20): suppose $\hat{\omega}=f \mathrm{dy}^{\mathrm{I}}+$ $g d y^{J} \wedge d x$ solves

$$
L_{1} \hat{\omega}=0
$$

(here $I$ and $J$ are fixed, $I=(1, J)$ ), and $\hat{\omega} \in H_{r}^{(2)}$ for some r. The comments following (3.21) indicate that $\hat{\omega} \in L^{2}\left(x^{-n} d x\right)$ and decreases rapidly at infinity. Then

$$
0=\int_{0}^{\infty}\left\langle L_{1} \hat{\omega}, \hat{\omega}\right\rangle x^{-n} d x=\int_{0}^{\infty}\left\langle L_{d} L_{\delta} \hat{\omega}+L_{\delta} L_{d} \hat{\omega}, \hat{\omega}\right\rangle x^{-n} d x
$$

$$
=\int_{0}^{\infty}\left(\left|L_{\delta} \hat{\omega}\right|^{2}+\left|L_{d} \hat{\omega}\right|^{2}\right) x^{-n} d x
$$

provided
a) $\int_{0}^{\infty} \frac{d}{d x}\left[x^{1-n}\left(x \frac{d g}{d x}-(n-k) g+(-1)^{k+1} i x|\eta| f\right) \bar{g}\right] d x=0$
and
b) $\int_{0}^{\infty} \frac{d}{d x}\left[x^{1-n}\left(x \frac{d f}{d x}-k f+(-1)^{k} i x|\eta| g\right) \bar{f}\right] d x=0$.

Now $L_{1}$ is a Fuchsian operator, the indicial roots of which coincide with those of the indicial operator of $P$. From (3.6) and a bit of calculation, if $k<\frac{n-1}{2}, f$ and $g$ have convergent expansions

$$
f=c_{0} x^{n-k-1}+\cdots+\log x \cdot\left(c_{j} x^{n-k+1}+\cdots\right)
$$

$$
\begin{equation*}
g=d_{0} x^{n-k}+\cdots+\log x \cdot\left(d_{0} x^{n-k}+\cdots\right) \tag{3.25}
\end{equation*}
$$

The logarithmic factors arise since the indicial roots differ by a positive integer. That these factors actually occur follows from an explicit computation of the first few terms of the Frobenius series. For later reference we
record that when $k>\frac{n+1}{2}$

$$
f=c_{0} x^{k}+\cdots+\log x \cdot\left(c_{0}^{\prime} x^{k}+\cdots\right)
$$

(3.26)

$$
g=d_{0} x^{k-1}+\cdots+\log x \cdot\left(d_{0} x^{k+1}+\cdots\right)
$$

To prove that the boundary contributions vanish, use these expansions near $x=0$ and the rapid decrease as $x \longrightarrow \infty$. Then, for example in $a)$, when $k<\frac{n-1}{2}$, the integrand is

$$
\frac{d}{d x}\left(a_{0} x^{n-2 k+1}+\cdots\right)
$$

while in b) it is

$$
\frac{d}{d x}\left(b_{0} x^{n-2 k-1}+\cdots\right)
$$

By the second fundamental theorem of calculus the boundary terms vanish as desired; similar reasoning holds when $k$ $\frac{n+1}{2}$.

We have so far shown that an $L^{2}$ solution of $N_{p} \hat{\omega}=0$ also satisfies

$$
\mathrm{L}_{\mathrm{d}} \hat{\omega}=0, \quad \mathrm{~L}_{\delta} \hat{\omega}=0
$$

In terms of the coordinate components $f$ and $g$ these
equations are

$$
\begin{gathered}
x \frac{d f}{d x}-k f+(-1)^{k} i x|\eta| g=0 \\
x \frac{d g}{d x}-(n-k) g+(-1)^{k+1} i x|\eta| f=0
\end{gathered}
$$

These equations may actually be uncoupled and solved. Thus, with a bit of work one sees that

$$
\begin{gathered}
x^{2} \frac{d^{2} f}{d x^{2}}-n x \frac{d f}{d x}+\left(k(n-k+1)-x^{2}|\eta|^{2}\right) f=0 \\
x^{2} \frac{d^{2} g}{d x^{2}}-n x \frac{d g}{d x}+\left((k+1)(n-k)-x^{2}|\eta|^{2}\right) g=0 .
\end{gathered}
$$

Once again $f$ and $g$ are Bessel functions, [19]:

$$
\mathrm{f}=\mathrm{c}_{1} \mathrm{x}^{\frac{\mathrm{n}+1}{2}} \mathrm{I}_{v_{2}}(\mathrm{x}|\eta|)+\mathrm{c}_{2} \mathrm{x}^{\frac{\mathrm{n}+1}{2}} \mathrm{~K}_{v_{2}}(\mathrm{x}|\eta|) \quad v_{2}=\left|\frac{\mathrm{n}-2 \mathrm{k}+1}{2}\right|
$$

(3.27)

$$
\mathrm{g}=\mathrm{d}_{1} \mathrm{x}^{\frac{\mathrm{n}+1}{2}} \mathrm{I}_{v_{1}}(\mathrm{x}|\eta|)+\mathrm{d}_{2} \mathrm{x}^{\frac{\mathrm{n}+1}{2}} \mathrm{~K}_{v_{1}}(\mathrm{x}|\eta|) \quad v_{1}=\left|\frac{\mathrm{n}-2 \mathrm{k}-1}{2}\right| .
$$

In order that $f, g$ lie in $L^{2}\left(x^{-n} d x\right)$ near zero, $c_{2}$ and $d_{2}$ must vanish, as in (3.13). But the solutions thus obtained increase exponentially, which is a contradiction. This proves Proposition (3.20).

Notice that the solutions $f$ and $g$ with $c_{2}=$ $d_{2}=0$ above are the factors of the logarithmic terms in (3.25) and (3.26). Finally we may complete the

Proof of (3.10): By virtue of (3.16) and (3.21) there is a kernel $G(x, \tilde{x}, \eta)$ for which

$$
\mathrm{N}_{\mathrm{P}}\left(\mathrm{x} \partial_{\mathrm{x}}, \mathrm{ix} \eta\right) \mathrm{G}(\mathrm{x}, \tilde{\mathrm{x}}, \eta)=\mathrm{x}^{\mathrm{n}} \delta(\mathrm{x}-\tilde{\mathrm{x}})
$$

and such that, for each $\eta$, the corresponding operator

$$
G: L^{2}\left(x^{-n} d x\right) \longrightarrow L^{2}\left(x^{-n} d x\right)
$$

is bounded. As of yet we may not assert this uniformly in $\eta$, which is what we require. Of course this uniformity is obvious for the 'uncoupled part' of $G$ (3.14); for the ' coupled part' we know little save that it is an analytic solution of $N_{P}$, of the form (3.25) for $x<\tilde{x}$ and (3.26) for $x>\tilde{x}$. For this extra information we use the dilation invariance of $\mathrm{N}_{\mathrm{P}}$.

Thus let $G(x, \tilde{x}, \tilde{\eta})$ be the solution above when $|\hat{\eta}|=$ 1. It is obviously smooth in $\hat{\eta}$. The solution $G$ to

$$
\mathrm{N}_{\mathrm{P}}\left(\mathrm{x} \partial_{\mathrm{x}}, \mathrm{ix}|\eta| \hat{\eta}\right) \mathrm{G}=|\eta|^{\mathrm{n}-1} \mathrm{x}^{\mathrm{n}} \delta(\mathrm{x}-\tilde{\mathrm{x}})
$$

obtained by replacing $x, \tilde{x}$ by $x|\eta|, \tilde{x}|\eta|$ must be

$$
\begin{equation*}
\mathrm{G}(\mathrm{x}, \tilde{\mathrm{x}}, \eta)=|\eta|^{1-\mathrm{n}} \mathrm{G}(\mathrm{x}|\eta|, \tilde{\mathrm{x}}|\eta|, \eta /|\eta|) \tag{3.28}
\end{equation*}
$$

where $\hat{\eta}=\eta /|\eta|$. The known part of the kernel (3.14) does indeed satisfy this.

The $L^{2}$ boundedness of $G$ as in (3.10) is now quite easy. Define the transformation

$$
D_{a}: f(x) \longrightarrow f_{a}(x)=a^{\frac{1-n}{2}} f(a x), \quad a \in \mathbb{R}^{+} .
$$

This correspondence is readily seen to be an $L^{2}\left(x^{-n} d x\right)$ isometry. Set $f a$ into the inequality

$$
\begin{equation*}
\int\left|\int \mathrm{G}(\mathrm{x}, \tilde{\mathrm{x}}, \hat{\eta}) \mathrm{f}(\tilde{\mathrm{x}}) \tilde{x}^{-\mathrm{n}} \mathrm{~d} \tilde{\mathrm{x}}\right|^{2} \mathrm{x}^{-\mathrm{n}} \mathrm{dx} \leq \mathrm{C} \int|\mathrm{f}|^{2} \mathrm{x}^{-\mathrm{n}} \mathrm{dx} \tag{3.29}
\end{equation*}
$$

where $C$ may be assumed independent of $\tilde{\eta} \in \mathrm{S}^{\mathrm{n}-2}$. The new right hand side, with $f a$ replacing $f$, is independent of a. On the other side, a change of variables $\tilde{x} \longmapsto a \tilde{x}$ results in a new inner integral on the left side

$$
\int G(x, \tilde{x}, \hat{\eta}) f_{a}(\tilde{x}) \tilde{x}^{-n} d \tilde{x}=\int a^{\frac{n-1}{2}} G\left(x, a^{-1} \tilde{x}, \hat{\eta}\right) f(\tilde{x}) \tilde{x}^{-n} d \tilde{x}
$$

Applying the dilation $\mathrm{D}_{\mathrm{a}^{-1}}$ transforms this to

$$
\int a^{n-1} G\left(a^{-1} x, a^{-1} \tilde{x}, \hat{\eta}\right) f(\tilde{x}) \tilde{x}^{-n} d \tilde{x}
$$

which retains the same $L^{2}\left(x^{-n} d x\right)$ norm. Thus, with $a=$ $|\eta|^{-1}$, the definition (3.28) implies
(3.30) $\int\left|\int G(x, \tilde{x}, \eta) f(\tilde{x}) \tilde{x}^{-n} d \tilde{x}\right|^{2} x^{-n} d x \leq C \int|f(x)|^{2} x^{-n} d x$
where $C$ is independent of $\eta \in \mathbb{R}^{n-1} \backslash\{0\}$. If $f$ is now allowed to depend on $\eta$ also, then by integrating both sides of (3.30) with respect to $\eta$ we obtain the required estimate of (3.10).

The kernel inverting the partial differential operator $N_{p}$ is

$$
\begin{equation*}
G(x, y, \tilde{x}, \tilde{y})=c_{n} \int e^{i(y-\tilde{y}) \cdot \eta_{G}}(x, \tilde{x}, \eta) d \eta \tag{3.31}
\end{equation*}
$$

It is bounded on $L^{2}\left(x^{-n} d x d y\right)$, hence is the unique Green function for the constant curvature Laplacian. It also has some isometry invariance, namely that of translation and dilation:

$$
\begin{gathered}
G(x, y+v, \tilde{x}, \tilde{y}+v)=G(x, y, \tilde{x}, \tilde{y}) \\
G(a z, a \tilde{z})=G(z, \tilde{z}), \quad z=(x, y), \quad \tilde{z}=(\tilde{x}, \tilde{y})
\end{gathered}
$$

the latter following from a change of variables in (3.31).
C. Structure of the Kernel

We now discuss the structure of the kernel $G$ of
(3.31), which is of the type discussed in the last chapter.
(3.32) Theorem: $G(z, \tilde{z}) \cdot \gamma \in K_{0}^{-2, \sigma, \tau}\left(\mathbb{R}_{+}^{n} \times_{0} \mathbb{R}_{+}^{n} ; \Gamma_{0}^{1 / 2} \otimes{ }^{o} \Lambda^{\mathrm{k}}\right)$ where

$$
\begin{aligned}
& \sigma=\left[\begin{array}{cc}
\mathrm{n}-\mathrm{k}-1 & \mathrm{n}-\mathrm{k}+1 \\
\mathrm{n}-\mathrm{k} & \mathrm{n}-\mathrm{k}
\end{array}\right], \quad \tau=\sigma^{\mathrm{t}}, \quad \text { for } \mathrm{k}<\frac{\mathrm{n}-1}{2}, \\
& \sigma=\left[\begin{array}{cc}
\mathrm{k} & \mathrm{k} \\
\mathrm{k}-1 & \mathrm{k}+1
\end{array}\right], \quad \tau=\sigma^{\mathrm{t}}, \quad \text { for } \mathrm{k}>\frac{\mathrm{n}-1}{2}, \quad \text { and } \quad \gamma
\end{aligned}
$$

is the half-density of (2.28).

As always, the proof is rather involved. First we prove uniform symbol estimates in the regions $\mathbf{x} \leq \tilde{x}, \tilde{x} \leq x$. This implies the requisite regularity in the interior and, together with properties of the indicial operator, yields the correct estimates at the boundary.

We commence by examining the 'uncoupled' part of the kernel:
(3.33) Lemma: The kernels $G_{1}$ and $G_{2}$ of (3.14), which furnish inverses for the operators $P_{1}$ and $P_{2}$, satisfy the estimates, for $|\eta| \geq 1$.

$$
\mid \partial_{\mathrm{x}}^{\mathrm{i}} \underset{\mathrm{x}}{\mathrm{j}} \partial_{\eta}^{\alpha} \mathrm{G}_{\mathrm{i}}\left(\mathrm{x}, \underset{\mathrm{x}, \eta)}{\sim} \mid \leq \mathrm{C}_{\epsilon, \mathrm{i}, \mathrm{j}, \alpha}(1+|\eta|)^{-1-|\alpha|+\mathrm{i}+\mathrm{j}_{\mathrm{e}}-|\mathrm{x}-\tilde{\mathrm{x}}||\eta|}\right.
$$

uniformly in $\epsilon \leq \mathbf{x} \leq \tilde{x} \leq \epsilon^{-1}, \epsilon \leq \tilde{x} \leq \mathbf{x} \leq \epsilon^{-1}$ for any $\epsilon>0, i=1,2$. In particular, neglecting the $(x, \tilde{x})$ dependence, they are symbols of order -1, uniformly in these regions of the $(x, \tilde{x})$ plane. Finally, each $G_{i}$ has a classical expansion as $\eta \longrightarrow 0$ of lowest homogeneity $|\eta|^{0}$.

Proof: Consider only $G_{1}$ in the region $\mathbf{x} \leq \tilde{x}$; the other cases are quite similar. Then

$$
\mathrm{G}_{1}(\mathrm{x}, \tilde{\mathrm{x}}, \eta)=\left(\tilde{x}^{\left.\left.\frac{\mathrm{x}}{}\right)^{\frac{\mathrm{n}-1}{2}} \mathrm{I}_{v}(\mathrm{x}|\eta|) \mathrm{K}_{v}(\tilde{\mathrm{x}}|\eta|), ~\right)}\right.
$$

where $v=(n-2 k-1) / 2$ for $k<(n-1) / 2$, which we also assume. Both $I_{v}$ and $K_{v}$ have asymptotic expansions for large values of their arguments, the first terms of which are given in (3.12). All derivatives of $I_{v}, K_{v}$ also have expansions which must equal the ones obtained by differentiating those for $I_{v}, K_{v}$. All these expansions are of step size one in descending powers of $x|\eta|$ or $\tilde{x}|\eta|$, hence it suffices to argue only for the first terms. Thus replace $G_{1}$ by

$$
\left(\underset{x x}{\tilde{x}}{ }^{\frac{\mathrm{n}-2}{2}} \mathrm{e}^{-(\tilde{x}-x)|\eta|}|\eta|^{-1}\right.
$$

as $|\eta| \longrightarrow \infty, \epsilon \leq \mathbf{x} \leq \tilde{\mathrm{x}}$. The estimates of the Lemma are quite easy to verify; the only point to be noted is that

$$
\partial_{\eta}^{\alpha} \mathrm{e}^{-(\tilde{\mathrm{x}}-\mathrm{x})|\eta|}=\sum_{|\beta| \leq|\alpha|} \mathrm{p}_{\beta}(\eta, \mathrm{x}-\tilde{\mathrm{x}}) \mathrm{e}^{-(\tilde{\mathrm{x}}-\mathrm{x})|\eta|}
$$

where each $p_{\beta}=(\tilde{x}-x)|\beta|$. $\{$ a term homogeneous in $\eta$ of degree $|\beta|-|\alpha|\}$. Since

$$
|\eta|^{|\beta|-|\alpha|} \underset{(\tilde{\mathrm{x}}-\mathrm{x})}{|\beta|_{\mathbf{e}^{-(\tilde{\mathrm{x}}-\mathrm{x})}|\eta|} \leq \mathrm{C}_{\beta}(1+|\eta|)^{-|\alpha|}}
$$

this too has the correct growth.
As for the behaviour as $\eta \longrightarrow 0$, both $I_{v}$ and $K_{v}$ have convergent Frobenius series (the one for $K_{v}$ containing logarithms) so that from (3.13)

$$
G_{1}(x, \tilde{x}, \eta)=c_{0} x^{n-k-1} \tilde{x}^{k}+|\eta|\left(c_{1} x^{n-k_{x}^{\sim} x^{n}} c_{1} x^{n-k-1} \tilde{x}^{k+1}\right)+\cdots
$$

(with an additional $|\eta| \log |\eta|$ term when $k=\frac{n}{2}-1$ ). This proves all claims.

To extend this to the coupled' part of the kernel $G_{c}$, which inverts $L_{1}$, we must understand its structure better. This necessitates a closer examination of the solutions to $L_{1} u=0$. Altogether, this nullspace is four dimensional and two independent elements are identified in
(3.27). In one of these, both components are exponentially increasing -call this solution $v_{1}$; the other solution $u_{1}$ has both components exponentially decreasing. Let us record that when $|\eta|=1$

$$
u_{1}=\left[\begin{array}{l}
a_{0} x^{n-k+1}+\cdots \\
b_{0} x^{n-k}+\cdots
\end{array}\right], \quad u_{1} \sim\left[\begin{array}{l}
a_{0} x^{n / 2} e^{x}+\cdots \\
b_{0} x^{n / 2} e^{x_{+}} \ldots
\end{array}\right] \quad x \longrightarrow \infty
$$

(3.34)

$$
v_{1}=\left[\begin{array}{l}
c_{0} x^{k}+\cdots \\
d_{0} x^{k+1}+\cdots
\end{array}\right], \quad v_{1} \sim\left[\begin{array}{l}
c_{0} x^{n / 2} e^{-x}+\cdots \\
d_{0} x^{n / 2} e^{-x_{+}} \ldots
\end{array}\right] x \longrightarrow \infty
$$

Other solutions are not known explicitly, but for tunately their asymptotics may be derived as follows. A Frobenius series calculation posits two additional solustions of the form (also when $|\eta|=1$ )

$$
u_{1}=\left[\begin{array}{ll}
\bar{a}_{0} & x^{n-k-1}+\cdots \\
\bar{b}_{0} & x^{n-k}+\cdots
\end{array}\right]+A u_{1} \log x
$$

(3.35)

$$
v_{2}=\left[\begin{array}{ll}
\bar{c}_{0} & x^{k}+\cdots \\
\bar{d}_{0} & x^{k-1}+\ldots
\end{array}\right]+B v_{1} \log x
$$

In (3.34), (3.35) none of the coefficients $a_{0}, a_{0}, \bar{a}_{0}, A$, $b_{0}, b_{0}^{\prime}, \bar{b}_{0}, B, c_{0}, c_{0}^{\prime}, \bar{c}_{0}, d_{0}, d_{0}^{\prime}, \bar{d}_{0}$ vanish.

On the other hand, the singularity of $L_{1}$ at $x=\infty$ is not regular. One may find formal solutions which are (possibly divergent) series in descending powers of $x$ multiplied by an exponential. Thus a straightforward computation produces series

$$
x^{\frac{n}{2}-1} e^{x}\left[\begin{array}{l}
\alpha_{1} x+\alpha_{0}+\alpha_{-1} x^{-1}+\cdots \\
\beta_{1} x+\beta_{0}+\cdots
\end{array}\right]
$$

$$
x^{\frac{n}{2}-1} e^{x}\left[\begin{array}{l}
\alpha_{-1} x^{-1}+\alpha_{-2} x^{-2}+\cdots \\
\beta_{-1} x^{-1}+\cdots
\end{array}\right]
$$

$$
x^{\frac{n}{2}-1} e^{x}\left[\begin{array}{c}
\alpha_{1} x+\cdots \\
\beta_{1} x+\cdots
\end{array}\right], \quad x^{\frac{n}{2}-1} e^{x}\left[\begin{array}{l}
\alpha_{-1} x^{-1}+\cdots \\
\beta_{-1} x^{-1}+\cdots
\end{array}\right]
$$

The $\alpha_{i}, \beta_{i}$ are different of course in each of these solutions, and none of the leading coefficients vanish. Once these formal solutions have been found, the theory developed in Chapter 5 of [8] (for which the analyticity of the coefficients of $L_{1}$ is essential) implies the existence of actual solutions obeying each of these asymptotics as $x \longrightarrow \infty$. Indeed both $u_{1}$ and $v_{1}$ of (3.34) fit this patern-their series being the ones commencing in $x^{1}$. Thus for some choice of $\overline{\mathrm{a}}_{\mathrm{i}}, \overline{\mathrm{b}}_{\mathrm{i}}, \overline{\mathrm{c}}_{\mathrm{i}}, \overline{\mathrm{d}}_{\mathrm{i}}$ in (3.35)

$$
u_{2} \sim\left[\begin{array}{c}
a_{0} x^{\frac{n}{2}-2} e^{x}+\cdots \\
b_{0}^{\prime} x^{\frac{n}{2}-2} e^{x}+\cdots
\end{array}\right], \quad a_{0}^{j}, b_{j}^{j} \neq 0
$$

(3.36)

$$
v_{2} \sim\left[\begin{array}{c}
c_{0}^{\prime} x^{\frac{n}{2}-2} e^{-x}+\cdots \\
d_{j} x^{\frac{n}{2}-2} e^{-x}+\ldots
\end{array}\right], \quad c_{0}^{j}, d_{0}^{j} \neq 0
$$

The kernel $G_{c}(x, \tilde{x}, \hat{\eta}), \quad|\hat{\eta}|=1$, has a form similar to (3.14):

$$
\mathrm{G}_{\mathrm{c}}(\mathrm{x}, \tilde{\mathrm{x}}, \hat{\eta})=\mathrm{U}(\mathrm{x}) \mathrm{A}(\tilde{x}) \mathrm{H}(\tilde{x}-\mathrm{x})+\mathrm{V}(\mathrm{x}) \mathrm{B}(\tilde{\mathrm{x}}) H(\mathrm{x}-\tilde{\mathrm{x}})
$$

where $U, V$ are two by two matrices, the columns of which are linear combinations of $u_{1}, u_{2}$, respectively $v_{1}, v_{2}$, and $A, B$ are matrices chosen so that

$$
\mathrm{L}_{1} \mathrm{G}_{\mathrm{c}}=\mathrm{x}^{\mathrm{n}} \delta(\mathrm{x}-\tilde{\mathrm{x}}) \cdot \mathrm{I}
$$

But $L_{1}$ is self-adjoint with respect to $x^{-n} d x$ (both formally and by (3.24) as an operator), so that $G_{c}$ is too:

$$
G_{c}(x, \tilde{x})=G_{c}^{*}(\tilde{x}, x)
$$

Hence $A(\tilde{x})=V(\tilde{x})^{*}, \quad B(\tilde{x})=U(\tilde{x})^{*}$ and
(3.37)

$$
G_{c}=U(x) V(\tilde{x})^{*} H(\tilde{x}-x)+V(x) U(\tilde{x}) H(x-\tilde{x}) .
$$

The only obstruction to completing our task is the explicit identification of the columns of $U, V$ in terms of the $u_{i}, v_{i}$. However, the condition that $L_{1} G_{c}=x^{n} \delta$ entails that

$$
U^{\prime}(x) V(x)^{*}-V^{\prime}(x) U(x)^{*}=-\frac{1}{2} x^{n-2} I
$$

Of necessity then, from (3.34) and (3.36)

$$
U(x)=\left[\begin{array}{ll}
\mathrm{au}_{2} & \mathrm{bu}_{1}
\end{array}\right], \quad \mathrm{V}(\mathrm{x})=\left[\begin{array}{ll}
\mathrm{cv}_{1} & \mathrm{~d} \mathrm{v}_{2} \tag{3.38}
\end{array}\right]
$$

for some $a, b, c, d$, none of which may vanish. Since

$$
\mathrm{G}_{\mathrm{c}}(\mathrm{x}, \tilde{\mathrm{x}}, \eta)=|\eta|^{1-\mathrm{n}_{\mathrm{G}}}(\mathrm{x}|\eta|, \tilde{\mathrm{x}}|\eta|, \eta /|\eta|)
$$

a verbatim repeat of the proof of (3.33) shows that
(3.39) Lemma: $G_{c}$ satisfies the estimates, for $|\eta| \geq 1$,

$$
\mid \partial_{\mathrm{x}}^{\mathrm{i}}{\underset{\sim}{\mathbf{x}}}_{j}^{\partial_{\eta}^{\alpha}} \mathrm{G}_{\mathrm{c}}(\mathrm{x}, \tilde{\mathrm{x}}, \eta) \mathrm{X} \leq \mathrm{C}_{\epsilon, \mathrm{i}, \mathrm{j}, \alpha}(1+|\eta|)^{-1-|\alpha|+\mathrm{i}+\mathrm{j}} \mathrm{e}^{-|\mathrm{x}-\tilde{\mathrm{x}}||\eta|}
$$ uniformly in $\epsilon \leq \mathbf{x} \leq \tilde{x} \leq \epsilon^{-1}, \epsilon \leq \tilde{x} \leq x \leq \epsilon^{-1}$. Furthermore, each term of $G_{c}$ has a classical expansion (in fact

a convergent series) as $\eta \longrightarrow 0$, uniform in $x, \tilde{x} \leq C$ and of lowest homogeneity $|\eta|^{0}$.

We now have sufficient information to study the kernel $G(z, \tilde{z})$ of (3.31).
(3.40) Lemma: $G(z, \tilde{z})$ is $C^{\infty}$ when $z \neq \tilde{z}$ and $x, \tilde{x}>0$. The singularity along the diagonal is classical conormal, with an expansion involving (quasi-) homogeneous terms of increasing degree in $z-\tilde{z}$; the first such term, up to a constant factor, is $|z-\tilde{z}|^{2-n}$.

Proof: By virtue of the uniform symbol estimates of the two previous lemmas, the usual arguments imply that $G(z, \tilde{z})$ is $C^{\infty}$ in the stated region. The singularity along the diagonal is the normal one for the inverse of an elliptic differential operator. From our vantage point, though, the symbolic behaviour of $G(x, \tilde{x}, \eta)$ away from $\eta=0$ shows that when $x=\tilde{x} G(z, \tilde{z})$ has an appropriate expansion in powers of $y-\tilde{y}$, the most singular of which has degree of homogeneity $-(n-1)-(-1)=2-n$. That the singularity behaves thusly in all directions along $\{z=\tilde{z}\}$ follows by invoking the rotation invariance of $G$. Alternately one might consider the Fourier transform of $G(x, x+t, \eta) \rho(t)$ with respect to $t$ (where $\rho$ is a cutoff function identically one near $t=0$ and compactly supported in,
say, $\left.\left(-\frac{1}{2} x, \frac{1}{2} x\right)\right)$. It is straightforward to show that this is symbolic in ( $\tau, \eta$ ) of degree -2 ( $\tau$ being dual to $t$ ) and polyhomogeneous.
(3.41) Proposition: $G(z, \tilde{z})$ is conormal on $\{x=0$, $\tilde{x}>0\},\{\tilde{x}=0, x>0\}$ and $\{\mathbf{x}=\tilde{x}=0, y \neq \tilde{y}\}$. Its components each have expansions on either of these codimension one faces, with leading terms determined by the two-bytwo matrices $\sigma$ for $\{\mathbf{x}=0\}$ and $\tau$ for $\{\tilde{\mathbf{x}}=0\}$ of (3.32).

Proof: We prove all assertions only for the inverse transform of the kernel $G_{1}$ of (3.14), and assuming $k<\frac{n-1}{2}$. The methods apply equally well to the full kernel, though the notation becomes more intricate.

Consider first the region near $x=0, \tilde{x} \geq \epsilon$. From (3.33) the integral

$$
\begin{aligned}
& \left(x \partial_{x}\right)^{\ell} \underset{\sim}{\underset{x}{j}}{\underset{y}{y}}_{\alpha}^{\partial_{\sim}^{\beta}} \underset{\sim}{\beta} \int e^{i(y-\tilde{y}) \cdot \eta_{G}}(x, \tilde{x}, \eta) d \eta \\
& =\int e^{i(y-\tilde{y}) \cdot \eta}\left(x \partial_{x}\right)^{\ell} \underset{\sim}{\underset{x}{j}(i \eta)^{\alpha}}(-i \eta)^{\beta}\left(x \tilde{x}^{\frac{n-1}{2}} I_{v}(|\eta|) K_{v}(\tilde{x}|\eta|) d \eta\right.
\end{aligned}
$$

converges absolutely for any $\ell, j, \alpha, \beta$, and uniformly in $x<\frac{1}{2} \epsilon, \epsilon \leq \tilde{x} \leq \epsilon^{-1},|y|,|\tilde{y}| \leq C ;$ indeed it decreases
exponentially there. Furthermore, it obviously has the same regularity as $x \longrightarrow 0$ independently of these various indices. This yields the conormality along $\{x=0, \tilde{x}>0\}$ since $\epsilon$ was arbitrary. Similar arguments apply to $\{\tilde{\mathrm{x}}=0, \mathrm{x}>0\}$.

For the corner, where both $x, \tilde{x} \sim 0$, it is convenient to observe that we may assume $|y-\tilde{y}|>\epsilon$, since the dilation invariance $G(a z, a \tilde{z})=G(z, \tilde{z})$ allows for all computations to be performed well away from $\{x=\tilde{x}=$ $y-\tilde{y}=0\}$. Now, assuming only that $x \leq \tilde{x}$, we may multiply both sides of the equation above by $(y-\tilde{y})^{\gamma}$ and replace $\underset{\tilde{x}}{\partial_{\tilde{x}}}$ by $\underset{\tilde{x}}{\tilde{x}} \partial_{\tilde{x}}$ to get

$$
\begin{aligned}
& \int e^{i(y-y) \cdot \eta}\left(-i \partial_{\eta}\right)^{\gamma} \\
& \quad\left[\left(x \partial_{x}\right)^{\ell} \tilde{x}_{\tilde{x} \partial_{\sim}^{x}}\right)^{j}(i \eta)^{\alpha}(-i \eta)^{\beta}\left(x \tilde{x}^{\frac{n-1}{2}} I_{v}(x|\eta|) K_{v}(\tilde{x}(|\eta|)] d \eta\right.
\end{aligned}
$$

as the new right hand side. (The integration by parts in $\eta$ is permissible since the whole expression is to be considered as an oscillatory integral.) Now, by choosing $\gamma$ so that $-1+\ell+j+|\alpha|+|\beta|-|\gamma|<-n$, the integral over $|\eta| \geq 1$ is convergent, whereas the integral over $|\eta| \leq 1$ can only produce a $C^{\infty}$ contribution.
Furthermore, the regularity of this expression is stable as
$x, \tilde{x} \longrightarrow 0$. Since $y-\tilde{y} \neq 0$, the conormality along the corner is now established.

As for the various expansions, notice that when $x \longrightarrow 0, \tilde{x} \geq \epsilon$, the first factor in

$$
\int I_{v}(x|\eta|) K_{v}(\tilde{x}|\eta|)\left(x \tilde{x}^{\frac{n-1}{2}} \mathrm{e}^{\mathrm{i}(\mathrm{y}-\tilde{\mathrm{y}}) \cdot \eta} \mathrm{d} \eta\right.
$$

may be replaced by its Frobenius series. Again, provided $x<\frac{1}{2} \epsilon$, each term of this (convergent) sum in the integrand decreases exponentially with $\eta$, hence we have a series

$$
\sum_{j=0}^{\infty} x^{-n-1+j} c_{j}(\tilde{x}, y-\tilde{y})
$$

furnishing the asymptotic behaviour as $x \longrightarrow 0$. To obtain this expansion uniformly down to $\tilde{x}=0$, but still assuming $x \leq \tilde{x}$ of course, we may instead replace $I_{v}(x|\eta|)$ by the first $N$ terms of its Frobenius series plus an exponentially increasing remainder. Each term in the resulting integrand exhibits symbolic behaviour, and the first $N$ summands still decrease exponentially in $\eta$. These provide the first $N$ terms of the asymptotic expansion; the remainder term is indeed lower order as may be seen by multiplying it by $(y-\tilde{y})^{\gamma}$ for $|\gamma|$ sufficiently
large and integrating by parts so that the resulting integral converges.

Applying the arguments of the last paragraph to the full kernel yields the correct expansions in all components, for these all are derived from (3.34), (3.35).

Proof of Theorem (3.32): The lift of $G(z, \tilde{z})$ to $\mathbb{R}_{+}^{n} \times_{0} \mathbb{R}_{+}^{n}$ has the correct behaviour along $\Delta \iota_{0}, T$, and $B$, as demonstrated in (3.40) and (3.41). Indeed the only point not yet discussed is the behaviour near the front face $F$. However this is trivial since the dilation invariance of G is equivalent to the independence of the lifted kernel on R. Hence not only is $G$ smooth on $F$, except at $\partial \Delta \iota_{0}$, but the expansions at $\Delta \iota_{0}, T, B$ all continue uniformly down to $F$.

Let us now untangle the various definitions and identify both the operator induced by the Laplacian on ${ }^{o} \Lambda^{k} \otimes$ $\Gamma_{0}^{1 / 2}$ and its inverse. First recall that the operator $P$, for which $G$ is an inverse, satisfies

$$
\Delta\left(x^{-k} \omega\right)=x^{-k}(P \omega) \Leftrightarrow P=x^{k} \Delta x^{-k}
$$

G then is the kernel induced by the standard Green function $G_{\Delta} \in \Gamma \operatorname{Hom}\left(\Lambda^{k}, \Lambda^{k}\right)$. In fact, since $\Lambda^{k}$ and ${ }^{\mathrm{o}} \Lambda^{\mathrm{k}}$ are canonically identified over the interior we may let $G_{\Delta}$ act on ${ }^{o} \Lambda^{k}$ :

$$
\begin{gathered}
\frac{\omega(x, y)}{x^{k}} \longrightarrow \int G_{\Delta}(x, y, \tilde{x}, \tilde{y}) \tilde{x}^{-k} \omega(\tilde{x}, \tilde{y}) \tilde{x}^{-n} d \tilde{x} d \tilde{y} \\
\quad=x^{-k} \int G(x, y, \tilde{x}, \tilde{y}) \omega(\tilde{x}, \tilde{y}) \tilde{x}^{-n} d \tilde{x} d \tilde{y}
\end{gathered}
$$

i.e.

$$
G_{\Delta}(x, y, \tilde{x}, \tilde{y})=x^{-k_{G}}(x, y, \tilde{x}, \tilde{y}) \tilde{x}^{\tilde{k}}
$$

(3.42)

$$
\Delta G_{\Delta}=x^{n} \delta(x-\tilde{x}) \delta(y-\tilde{y}) \cdot I .
$$

To incorporate half-densities, use $\mu$ and $\gamma$ of (2.28) and define $\tilde{P}$ by

$$
\Delta\left(\mathrm{x}^{-\mathrm{k}} \omega \cdot \mu\right)=\mathrm{x}^{-\mathrm{k}}(\tilde{\mathrm{P}} \omega) \cdot \mu \Leftrightarrow \tilde{\mathrm{P}}=\mathrm{x}^{\mathrm{k}+\mathrm{n} / 2} \Delta \mathrm{x}^{-\mathrm{k}-\mathrm{n} / 2}=\mathrm{x}^{\mathrm{n} / 2} \mathrm{Px}^{-\mathrm{n} / 2}
$$

It is of course the operator induced by $\Delta$ on ${ }^{0} \Omega^{k} \otimes \Gamma_{0}^{1 / 2}$. The corresponding kernel $\tilde{G}_{\Delta} \in \Gamma \operatorname{Hom}\left({ }^{\mathrm{o}} \Lambda^{\mathrm{k}},{ }^{\mathrm{o}} \Lambda^{\mathrm{k}}\right)$ must satisfy

$$
\Delta \cdot\left(\tilde{G}_{\Delta} \cdot r\right)=\left(x^{n / 2} \Delta x^{-n / 2} \tilde{G}_{\Delta}\right) \cdot r=x^{n} \delta(x-\tilde{x}) \delta(y-\tilde{y}) \cdot I
$$

-which is the kernel of the identity with respect to $\Gamma_{0}^{1 / 2}$. From (3.42) one may check that

$$
\begin{equation*}
\tilde{\mathrm{G}}_{\Delta}=\mathrm{x}^{\mathrm{n} / 2}{ }_{G_{\Delta}} \tilde{\mathrm{x}}^{-\mathrm{n} / 2}=\mathrm{x}^{-\mathrm{k}+\mathrm{n} / 2} \tilde{G x}^{\tilde{x}^{\mathrm{k}-\mathrm{n} / 2}} . \tag{3.43}
\end{equation*}
$$

The kernel $G_{\Delta}$ transfers to $B^{n}$ and is invariant under all hyperbolic isometries there. The preceding discussion may be carried through on the ball, the only difference being that $x$ must be replaced by $1-|w|^{2}$ (w is the Euclidean coordinate on $B^{n}$ ). The appropriate kernel inverting the Laplacian acting on. ${ }^{o} \Lambda^{k}\left(B^{n}\right) \otimes$ $\Gamma_{0}^{1 / 2}\left(\mathrm{~B}^{\mathrm{n}}\right)$ is

$$
\begin{equation*}
\tilde{G}_{\Delta}=\left(1-|w|^{2}\right)^{n / 2} G_{\Delta}(w, \tilde{w})\left(1-|\tilde{w}|^{2}\right)^{-n / 2} . \tag{3.44}
\end{equation*}
$$

Suppose now that $\sigma$ is any two-by-two matrix-the components of which regulate decay rates of the tangential and normal components of a section of $\operatorname{Hom}\left({ }^{0} \Lambda^{k}, \Lambda^{\mathrm{o}}{ }^{\mathrm{k}}\right)$-and define the new matrix

$$
\sigma+q=\left(\sigma_{i j}+q\right), \quad i, j=t, n, \quad q \in \mathbb{R}
$$

In summary, we have proved
(3.45) Corollary:

$$
\tilde{\mathrm{G}}_{\Delta}(\mathrm{w}, \tilde{\mathrm{w}}) \cdot r \in \mathrm{~K}_{0}^{-2, \sigma-\mathrm{k}+\mathrm{n} / 2, \tau+\mathrm{k}-\mathrm{n} / 2}\left(\mathrm{~B}^{\mathrm{n}} ;{ }^{\mathrm{o}} \Lambda^{\mathrm{k}} \otimes \Gamma_{0}^{1 / 2}\right)
$$

$\sigma, \tau$ as in (3.32). Furthermore

$$
\Delta\left(\tilde{G}_{\Delta} \cdot \gamma\right)=\left(1-|w|^{2}\right)^{n} \delta(w-\tilde{w}) \cdot \gamma \cdot I .
$$

D. Refined Mapping Properties

In this section we further refine the mapping properties of the Laplacian and its inverse. Two results are proved; the first, a straightforward extension of Proposition (2.55), is valid more generally, but its statement and proof are simpler in the present circumstances. Let us define a form valued version of the space (2.54).

$$
\begin{align*}
& \omega \in \alpha_{\mathrm{phg}}^{\alpha}\left(\mathrm{M},{ }^{\mathrm{o}} \Lambda^{\mathrm{k}}\right), \quad \alpha=\left(\alpha_{\mathrm{t}}, \alpha_{\mathrm{n}}\right) \\
& \Leftrightarrow \omega=\omega_{t}+\omega_{n} \in{ }^{o} \Omega_{t}^{k} \oplus{ }^{o} \Omega_{n}^{k} \text { near } \partial M  \tag{3.46}\\
& \omega_{i} \in{\underset{p h g}{ }}_{\alpha_{i}}(M) \otimes{ }^{o} \Omega_{i}^{k}, \quad i=t, n .
\end{align*}
$$

The decomposition ${ }^{o} \Lambda^{k}={ }^{o} \Lambda_{t}^{k} \oplus{ }^{o} \Lambda_{n}^{k}$ near $\quad \partial M$ is defined following the discussion in Chapter 2, Section $C$, and the spaces above are invariant under coordinate changes only when $\left|\alpha_{t}-\alpha_{n}\right| \leq 1$.
(3.47) Proposition: Let $f \in \mathscr{A l}_{\mathrm{phg}}^{\alpha+q}\left(\mathrm{~B}^{\mathrm{n}},{ }^{\mathrm{o}} \Lambda^{\mathrm{k}}\right)$, where $\alpha=$ $(\mathrm{n}-\mathrm{k}-1, \mathrm{n}-\mathrm{k})$ if $\mathrm{k}<(\mathrm{n}-1) / 2, \quad \alpha=(\mathrm{k}, \mathrm{k}-1)$ if $\mathrm{k}>(\mathrm{n}+1) / 2$ and $\alpha+q=\left(\alpha_{i}+q\right), q \in \mathbb{Z}^{+}$, in either case. Then there is a unique $\omega \in A_{p h g}^{\alpha}\left(B^{n},{ }^{o} \Lambda^{k}\right)$ such that

$$
\Delta \omega=f .
$$

Proof: Let $\rho$ be a defining function for $\partial \mathrm{B}^{\mathrm{n}}, \quad \rho \geq 0$, and $\theta$ the variable in $\partial B^{n}=S^{n-1}$. Now, the tangential and normal components of $f$ have expansions

$$
\begin{gathered}
f_{i} \sim \sum_{j=0}^{\infty} f_{i j} \quad i=t, n \\
f_{i j}=\sum_{\ell=0}^{N_{i j}} C_{i j \ell}(\theta) \rho^{\alpha_{i}+q+j}(\log \rho)^{\ell}
\end{gathered}
$$

near the boundary. Inasmuch as

$$
\Delta-I_{\Delta}: A^{\alpha} \longrightarrow A^{\alpha+1}
$$

for any $\alpha$, we proceed by first choosing $\omega_{i o}$, $i=t, n$ so that

$$
\begin{gathered}
\mathrm{I}_{\Delta} \omega_{\text {io }}=\mathrm{f}_{\mathrm{io}} \\
\omega_{\text {io }}=\sum_{\ell=0}^{\mathrm{N}_{\text {io }}} \mathrm{c}_{\text {io } \ell}(\theta) \rho^{\alpha_{i}+\mathrm{q}}(\log \rho)^{\ell} .
\end{gathered}
$$

Then, assuming $\omega_{i m}$ have been chosen, $m<j, p i c k \omega_{i j}$ so that

$$
\begin{gathered}
I_{\Delta}^{\omega_{i j}}=f_{i j}^{\prime} \\
\omega_{i j}=\sum_{\ell=0}^{N_{i j}} c_{i j \ell}^{\prime}(\theta) \rho^{\alpha_{i}+q+j}(\log \rho)^{\ell} .
\end{gathered}
$$

The $f_{i j}$ are quasihomogeneous terms of the correct degree which depend only on $f_{i j}$ and $\left(\Delta-I_{\Delta}\right) \omega_{t m},\left(\Delta-I_{\Delta}\right) \omega_{n m}$, $m<j$. All these equations are solvable by finite series methods. Note only that uniqueness holds at each step since each exponent $\alpha_{i}+q+j$ is greater than the largest of the indicial roots of $I_{P}$ - the operator induced by $I_{\Delta}$ on ${ }^{o} \Lambda^{k}$. For the same reason the highest power of the logarithm at each stage is never increased. If $\omega_{0} \sim \sum \omega_{i j}$ then $\Delta \omega_{0}-f \simeq g \in \dot{C}^{\infty}\left(B^{n}\right) \otimes{ }^{o} \Lambda^{k}$. (2.55) and (3.45) together guarantee the existence of a $\omega_{1}$ such that

$$
\Delta \omega_{1}=\mathrm{g}
$$

Then $\Delta\left(\omega_{0}{ }^{-\omega_{1}}\right)=f, \quad$ as desired.

The other solvability result is of the same type as (3.47), but with somewhat more singular right hand sides. Recall from (2.17) that it is useful to think of the quarter-sphere $S_{++}^{n}$ as the ball $B^{n}$ blown up around a point $w_{0}$ in its boundary. Analogous to the spaces (3.46)
are

$$
\begin{equation*}
\alpha_{\mathrm{phg}}^{\alpha, \beta}\left(\mathrm{S}_{++}^{\mathrm{n}}{ }^{\mathrm{o}} \Lambda^{\mathrm{k}}\right) \longleftrightarrow A_{\mathrm{phg}}^{\alpha, \beta}\left(\mathrm{B}^{\mathrm{n}},{ }^{\mathrm{o}} \Lambda^{\mathrm{k}}\right) \tag{3.48}
\end{equation*}
$$

$\alpha=\left(\alpha_{t}, \alpha_{n}\right), \beta=\left(\beta_{t}, \beta_{n}\right)$. An element is a section of ${ }^{o} \Lambda^{k}$ which is conormal, with classical expansion, at the top and bottom edges of the quarter-sphere. The two-vectors $\alpha, \beta$ determine the most singular terms in these expansions; as in (3.46) we require $\left|\alpha_{t}-\alpha_{n}\right|,\left|\beta_{t}-\beta_{n}\right|=0$ or 1. The analogous spaces over $B^{n}$ are those with expansions at $\partial B^{n}-\left\{w_{0}\right\}$ and at $\left\{W_{0}\right\}$. In particular, using the defining function $r(w)=\left|w-w_{0}\right|$, the tangential component of an element has an expansion

$$
\begin{equation*}
\mathbf{f}_{\mathrm{t}} \sim \sum_{\mathrm{i}=0}^{\infty} \sum_{\ell=0}^{\mathrm{N}_{\mathrm{i}}} \phi_{i \ell^{r}}^{\beta_{\mathrm{t}}+\mathrm{i}}(\log \mathrm{r})^{\ell} \quad \mathrm{r} \longrightarrow 0 \tag{3.49}
\end{equation*}
$$

and similarly for the normal component. The $\phi_{i} \ell$ here are functions of $\theta \in \mathrm{S}_{+}^{\mathrm{n}-1}$, the sphere in the half tangent space to $B^{n}$ at $w_{0}$. We may now state the
(3.50) Theorem: For any $f \in A_{p h g}^{\alpha+q, \beta}\left(B^{n},{ }^{o} \Lambda^{k}\right), \quad q \in \mathbb{Z}^{+}$, there is always a solution

$$
\omega \in A_{\mathrm{phg}}^{\alpha, \beta}\left(\mathrm{B}^{\mathrm{n}},{ }^{o} \Lambda^{\mathrm{k}}\right)
$$

to $\Delta \omega=f$, so long as the entries of $\beta$ are integers and $\alpha$ is the 2 -vector of (3.47).

The typically arduous proof involves removing the additional singularity at $w_{0}$ so as to reduce the problem to (3.47). This step is essentially local, so we replace $B^{n}$ by $\mathbb{R}_{+}^{n}$ with $w_{0}$ corresponding to 0 . In addition, replace $\Delta$ by $N_{P}$ and regard $\omega$ and $f$ as sections of $\Lambda^{\mathrm{k}}$ rather than ${ }^{o} \Lambda^{k}$. The principle tool is the Mellin transform, taken in the radial variable $r$, where $z \in \mathbb{R}_{+}^{n}$ has polar coordinates $r \in \mathbb{R}^{+}, \theta \in S_{+}^{n-1}$. This is the transform

$$
u(r, \theta) \longrightarrow u_{M}(\zeta, \theta)=\int r^{-\zeta-1} u(r, \theta) d r
$$

$$
\begin{equation*}
u(r, \theta)=(2 \pi)^{-1} \int r^{\zeta} u_{M}(\zeta, \theta) d \zeta . \tag{3.51}
\end{equation*}
$$

In general, $u_{M}$ is only defined (and analytic) in a halfspace $\operatorname{Re} 5<a$, and so the integral defining the inverse transform is taken over a line $\operatorname{Re} \zeta=c, c<a \quad$ fixed. (For all facts quoted here regarding the transform, see [23].) The constant a which limits the domain of ${ }^{u_{M}}$ is closely related to the decay of $u$ as $r \longrightarrow 0$. If $u=$ 0 for large values of $r$ and vanishes at definite rate at $r=0$, then $u_{M}$ is rapidly decreasing as $|\operatorname{Im} \zeta| \longrightarrow \infty, \operatorname{Re} \zeta \quad$ fixed.

The most important property of $u_{M}$ for $u s$ is that if $u(r, \theta)$ has an expansion as in (3.49), and vanishes for $r$ large, then $u_{M}$ continues meromorphically to the whole $\zeta$-plane, with poles only at $\zeta \in \beta_{0}+\mathbb{Z}^{+}$(if $\beta_{0}$ is the most singular exponent in the expansion). Furthermore, the order of the pole at $\beta_{0}+i$ is $N_{i}+1$. Conversely, if $\mathrm{v}(\zeta, \theta)$ is meromorphic in $\zeta$ with real poles and decreases rapidly on each line $\operatorname{Re} \zeta=$ constant, then

$$
u(r, \theta)=\int_{\operatorname{Re} \zeta=c} r^{\zeta} v(\zeta, \theta) d \zeta
$$

has an expansion as in (3.49) with the exponent $s$ in each term $r^{s}(\log r)^{l}$ the location of a pole of $v$. Only those $s$ occur for which Res $>$ c, and the order of the pole at $\zeta=s$ is one greater than the highest power of the logarithm multiplying $r^{s}$.

Define the operator $L_{\xi}$ by

$$
\left(N_{P} u\right)_{M}=L_{\zeta} u_{M}(\zeta, \theta)
$$

Clearly

$$
L_{\zeta}=|z|^{-\zeta}{N_{P}}|z|^{\zeta}
$$

Since $N_{P}$ is invariant under the homothety $z \longrightarrow a z$, this conjugate is also. Therefore it is differential only
in the $\theta$ variable; in $\zeta$ it is purely algebraic.
Furthermore, as an operator on $S_{+}^{n-1}=\left\{\left(\theta_{1}, \cdots, \theta_{n}\right)\right.$ : $\left.|\theta|=1, \theta_{n} \geq 0\right\}$ it is in $D i f f_{0}^{2}$ With respect to the boundary $\theta_{n}=0$ of course. A short calculation shows that its indicial operator is one we know:

$$
\begin{equation*}
I\left(L_{\zeta}\right)=I_{P} \tag{3.52}
\end{equation*}
$$

The present goal is to invert $L_{S}$, so as to solve the equation $L_{\zeta}{ }^{\omega} M=f_{M}$. Naturally, $L_{\zeta}$ is not invertible for every value of $\zeta$. Its inverse $G_{\zeta}$ is meromorphic, and its poles contribute to the expansion of $\omega_{M}$. Let $\rho(r) \epsilon$ $C_{0}^{\infty}(\mathbb{R})$ equal one for $|r| \leq 1 / 2$, and define, at least formally, the inverse for $L_{S}$

$$
\left.G_{\zeta}^{f}=\lim _{t \rightarrow 0} M_{t}^{*}\left\{r^{-\zeta_{G}\left(r^{\prime}\right.} \rho(r) f(\theta)\right)\right\}
$$

$f \in \dot{C}^{\infty}\left(S_{+}^{n-1}\right) \otimes{ }^{o} \Lambda^{k}, \quad G$ the inverse of $N_{P}$, and $M_{t}$ the dilation by $t$. Substituting the integral for $G$ in this formula we get

$$
\begin{equation*}
G_{\zeta} f(\theta)=\int_{S_{+}^{n-1}} H_{\zeta}(\theta, \tilde{\theta}) f(\tilde{\theta})\left(\tilde{\theta}_{n}\right)^{-(n-1)} d \tilde{\theta} \tag{3.53}
\end{equation*}
$$

where

$$
\begin{equation*}
\mathrm{H}_{\zeta}(\theta, \tilde{\theta})=\int_{0}^{\infty} \mathrm{G}(\theta, \tau \tilde{\theta}) \tau^{\zeta} \frac{\mathrm{d} \tau}{\tau \tilde{\theta}_{\mathrm{n}}} \tag{3.54}
\end{equation*}
$$

For any value $\zeta$ at which this last integral converges, $G_{\zeta}$ then serves as an actual inverse for $L_{S}$. In fact, but for an insignificant difference, $G_{\zeta}$ is actually a $V_{0}$-pseudodifferential operator. Let us now examine the values of $\zeta$ where (3.54) converges. Fix $\theta \neq \tilde{\theta}$ and notice that by the dilation invariance of $G, G(\theta, \tau \tilde{\theta})=$ $G\left(\tau^{-1} \theta, \tilde{\theta}\right)$. The components of $G$ decays at various rates as $\tau \longrightarrow 0$ and $\tau \longrightarrow \infty$. However, the extreme case is $G_{t}$; all other components exhibit 'more favourable' rates of decay. Suppose $k<\frac{n-1}{2}$; then

$$
\begin{gathered}
\mathrm{G}_{\mathrm{tt}}(\theta, \tau \tilde{\theta}) \sim \tau^{\mathrm{n}-\mathrm{k}-1,} \quad \tau \longrightarrow 0 \\
\mathrm{G}_{\mathrm{tt}}\left(\tau^{-1} \theta, \tilde{\theta}\right) \\
\sim \tau^{-(\mathrm{n}-\mathrm{k}-1)}, \quad \tau \longrightarrow \infty .
\end{gathered}
$$

Hence (3.54) converges when $-(n-k-1)<R e \zeta<n-k-1$.
(3.55) Lemma: $H_{\zeta}$ extends meromorphically to the whole $\zeta$-plane with poles at $\pm \zeta \in \mathrm{n}-\mathrm{k}-1+\mathbb{Z}^{+}$if $\mathrm{k}<\frac{\mathrm{n}-1}{2}$ and $\pm \zeta \epsilon k-1+\mathbb{Z}^{+}$if $k>\frac{n-1}{2}$. These poles are of at most order two; the residues at $\zeta=\mathrm{N}$ of $\mathrm{H}_{\zeta}$ and $(\zeta-\mathrm{N}) \mathrm{H}_{\zeta}$ are kernels of finite rank.

Proof: Choose a partition of unity $1=\phi_{1}(\tau)+\phi_{2}(\tau)+$ $\phi_{3}(\tau)$, where $\phi_{3}(\mathrm{t})=\phi_{1}\left(\tau^{-1}\right), \quad \phi_{1}$ is supported in $\tau \leq 1 / 2$, and $\phi_{2}$ in $(1 / 4,4)$. Write $H_{\zeta}=H_{\zeta}^{1}+H_{\zeta}^{2}+H_{\zeta}^{3}$ by letting $H_{\zeta}^{i}$ be defined by the integral (3.54) with $\phi_{i}(\tau)$ inserted into the integrand. $H_{\zeta}^{2}$ is entire in $\zeta$ for $\theta \neq \tilde{\theta}$, which singularity we discuss later. By a change of variables

$$
\begin{aligned}
\mathrm{H}_{\zeta}^{1}(\theta, \tilde{\theta}) & =\int \phi_{1}(\tau) \mathrm{G}(\theta, \tau \tilde{\theta}) \tau \zeta \frac{\mathrm{d} \tau}{\tau \tilde{\theta}_{\mathrm{n}}} \\
\mathrm{H}_{\zeta}^{3}(\theta, \tilde{\theta}) & =\int \phi_{1}(\tau) \mathrm{G}(\tau \theta, \tilde{\theta}) \tau-\zeta \frac{\mathrm{d} \tau}{\tau \tilde{\theta}_{\mathrm{n}}} .
\end{aligned}
$$

Replace $\tau^{\lambda}$ by $\tau_{+}^{\lambda}, \lambda= \pm \zeta$, in each of these expiression. The family $\tau_{+}^{\lambda}$ of homogeneous distributions is well known to extend meromorphically to $\mathbb{C}$ with simple poles at $\lambda \in-\mathbb{Z}^{+}$; the residue at $\lambda=-N$ is $(-1)^{N-1} \cdot \delta^{(N-1)} /(N-1)!$.

Now, replace $G(\theta, \tau \tilde{\theta})$ and $G(\tau \theta, \tilde{\theta})$ by their expansion as $\tau \longrightarrow 0$. These expansions involve terms $\tau^{m} G_{m o}$, $\tau^{m} \log \tau G_{m 1}$, where the $G_{m \ell}$ are kernels in $(\theta, \tilde{\theta})$ which are polynomial in $\theta$ or $\tilde{\theta}$ for $\mathrm{G}(\tau \theta, \tilde{\theta})$ or $\mathrm{G}(\theta, \tau \tilde{\theta})$, respectively. The expansions commence with $\tau^{\mathrm{n}-\mathrm{k}-1}, \mathrm{k}<$ $\frac{n-1}{2}$, or $\tau^{k-1}, k>\frac{n+1}{2}$. Inasmuch as $\tau_{+}^{\lambda} \tau^{m}=\tau_{+}^{\lambda+m}$ for any $\lambda, m$, and $\tau_{+}^{\lambda} \log \tau=d_{\lambda}\left(\tau_{+}^{\lambda}\right)$, the integrals of each term of the expansions should be thought of as a pairing of
$\tau_{+}^{\zeta+m}$ or $\tau_{+}^{\zeta+m} \log \tau$ with $\phi_{1}(\tau)$; the coefficient here is a kernel of finite rank since it is polynomial in $\theta$ or $\tilde{\theta}$. Finally, since $\phi_{1}^{(i)}(0)=0, i>0$, only one or two terms will actually be singular. Specifically, at $\zeta=$ - ( $n-k-1+i)$, for example, $H_{\zeta}^{1}$ will have a pole resulting only from the two terms $\tau^{n-k+i} G_{i o}(\theta, \tilde{\theta}), \tau^{n-k-1+i} \log \tau$ $G_{i 1}(\theta, \tilde{\theta})$ in the expansion for $G(\theta, \tau \tilde{\theta})$, and $H_{S}^{3}$ is regular. Similarly, the pole at $\zeta=n-k-1+i$ occurs only in $H_{\zeta}^{3}$, and not in $H_{\zeta}^{1}$. The proof is complete.

Since $L_{\zeta} G_{\zeta}=\left(\theta_{n}\right)^{n-1} \delta(\theta-\tilde{\theta})$ in the region where (3.54) converges, it must continue to hold in the extended domain, away from the poles, by the uniqueness of analytic continuation. The kernel $H_{S}$ essentially lies in $K_{0}^{-2, \sigma, \tau-1}$ over $S_{+}^{n-1} \times_{0} S_{+}^{n-1}$, as we shall now demonstrate. First of all, $H_{S}$ is the pushforward of $G$ (with a few extra $\tau$ factors) under the submersion $(\theta, \tilde{\theta}, \tau) \longrightarrow$ $(\theta, \tilde{\theta})$. The restriction of this map to the submanifold $\{\theta$ $=\tilde{\theta}\}$ is still a submersion onto its image, and $G(\theta, \tau \tilde{\theta})$ is conormal along $\{\theta=\tilde{\theta}, \tau=1\}$, so $H_{S}$ is conormal along $\{\theta=\tilde{\theta}\}$. Next, $H_{\zeta}$ has the proper singularities along the top and bottom faces, as may be ascertained by setting the expansion for $G$ into (3.54). Finally, the one way in which $H_{S}$ differs from the $V_{0}$ kernels we have already studied is that it has a logarithmic singularity at the front face. The easiest way to check this is to
observe that the kernel $G_{\Delta}$ of (3.42) is actually a function of the hyperbolic distance between $z$ and $\tilde{z}$. This simplifies the calculations necessary to display this singular behaviour, which are then straightforward and much the same as those in [20], hence will be omitted.

Despite this one difference, $H_{S}$ still enjoys all the continuity properties discussed in Chapter 2. From this, and using (3.52), the proof of (3.47) may be repeated to show that

$$
L_{\zeta^{\omega}}^{\omega}=f \in \alpha_{p h g}^{\alpha+q}\left(S_{+}^{n-1}, \Lambda^{k}\right)
$$

has a unique solution $\omega \in A_{p h g}^{\alpha}\left(S_{+}^{n-1}, \Lambda^{k}\right)$ depending meromorphically on $\zeta$. Its poles are precisely those of ${ }^{H}$.

Proof of (3.50): Take $f \in \mathcal{A}_{\mathrm{phg}}^{\alpha+q, \beta}\left(B^{n},{ }^{\mathrm{o}} \Lambda^{k}\right)$ and transfer it to the upper half-space, with $w_{0}=0$. Also, multiply it by $x^{k}$ so as to get an ordinary k-form. Now, express its components in polar coordinates, and take their Mellin transform with respect to $r$. The new form, $f_{M}(\zeta, \theta)$, is meromorphic; its poles are at $\zeta \in \min \left(\beta_{t}, \beta_{n}\right)+\mathbb{Z}^{+}$, and it is still in $A_{\mathrm{phg}}^{\gamma+\mathrm{q}}$ in $\theta$. Let $\omega_{\mathrm{M}}(\zeta, \theta)=\mathrm{G}_{\zeta} \mathrm{f}_{\mathrm{M}}(\zeta, \theta)$. It too is meromorphic in $\zeta$, and lies in $\alpha_{\mathrm{phg}}^{\alpha}$ in $\theta$. By virtue of the hypothesis that the $\beta_{i}$ are integers, ${ }^{\omega} M$
has poles at $\pm \zeta \in \mathrm{n}-\mathrm{k}-1+\mathbb{Z}^{+}$(or $\mathrm{k}-1+\mathbb{Z}^{+}$) and $\min \left(\beta_{t}, \beta_{n}\right)+\mathbb{Z}^{+}$.

Define

$$
\omega_{1}(r, \theta)=(2 \pi)^{-1} \int r^{\zeta} \omega_{M}(\zeta, \theta) d \zeta
$$

where the integral is over $\operatorname{Re} \zeta=\min \left(\beta_{t}, \beta_{n}\right)-1 / 2$. Then

$$
\mathrm{N}_{\mathrm{P}} \omega_{1}-\mathbf{f}=\mathbf{g}
$$

vanishes to infinite order at $r=0$, and upon being transferred back to $B^{n}$, lies in $A_{p h g}^{\alpha+q}\left(B^{n}, \Lambda^{k}\right)$. Now apply (3.47) to find $\omega_{2}$ such that

$$
\mathrm{N}_{\mathrm{P}} \omega_{2}=\mathrm{g}
$$

$\omega_{1}-\omega_{2}$ is the desired solution: $N_{P}\left(\omega_{1}-\omega_{2}\right)=f$.

Chapter 4: Hodge Cohomology
A. The Parametric Construction

We now have all the analytic tools at our disposal necessary to construct a suitable parametrix for the Laplacian on $k$-forms for a conformably compact metric (1.2). As described in section $D$ of Chapter 2, the construction proceeds in three stages. Each of these involve inverting model operators for $\Delta_{g}$-the symbol, the normal operator, and the indicial operator at the diagonal, the front face, and the top face of $M \times{ }_{0} M$, respectively.

The first step uses the symbol calculus of Theorem (2.41), but is otherwise identical to the usual microlocal procedure. Thus we seek an operator

$$
\mathrm{E}_{0} \in \Psi_{0}^{-2}\left(\mathrm{M} ;{ }^{\mathrm{o}} \Lambda^{\mathrm{k}} \otimes \Gamma_{0}^{1 / 2}\right)
$$

the kernel of which is supported quite near $\Delta \iota_{0}$ in $M \times{ }_{0} M$, and of course such that

$$
\begin{equation*}
\Delta_{\mathrm{g}} \mathrm{E}_{0}=\mathrm{I}-\mathrm{Q}_{0}, \quad \mathrm{Q}_{0} \in \Psi_{0}^{-\infty} \tag{4.1}
\end{equation*}
$$

First choose $E_{0,0} \in \Psi_{0}^{-2}$ satisfying ${ }^{o} \sigma_{2}\left(\Delta_{g}\right)^{o} \sigma_{-2}\left(E_{0,0}\right)=$ 1, which is possible by the $V_{0}$ ellipticity of $\Delta_{g}$.

Then

$$
\Delta_{g} E_{0,0}-I=-Q_{0,1} \in \Psi_{0}^{-1}
$$

Similarly, by induction, we may find $E_{0, j} \in \Psi_{0}^{-2-j}$ such that

$$
\Delta_{g} E_{0, j}-Q_{0, j}=-Q_{0, j+1} \in \Psi_{0}^{-j-1}
$$

$j=1,2, \cdots$ Obviously it may be assumed that the kernel of each $E_{0, j}$, and thus $Q_{0, j}$, is supported near $\Delta \iota_{0}$. Now take

$$
E_{0} \sim \sum E_{0, j}
$$

and the first step is complete.
Next we seek an operator $E_{1}$ for which

$$
\begin{equation*}
\Delta_{g} E_{1}=Q_{0}-Q_{1} \tag{4.2}
\end{equation*}
$$

where the kernel of $Q_{1}$ vanishes to infinite order on the front face $F$. Although $\kappa\left(Q_{0}\right)$ is supported well away from all other boundaries, it is not possible to choose $\mathrm{E}_{1}$ and $Q_{1}$ with this property. In fact, we will find that
(4.3) $\quad E_{1} \in \Psi_{0}^{-\infty}, \sigma^{\prime}, \tau^{\prime}, \quad Q_{1} \in R^{\infty} \Psi_{0}^{-\infty}, \sigma^{\prime}, \tau^{\prime}$
where $\sigma^{\prime}=\sigma-\mathrm{k}+\mathrm{n} / 2, \tau^{\prime}=\tau+\mathrm{k}-\mathrm{n} / 2, \sigma$ and $\tau$ the matrices of (3.32). We note that it is precisely at this step where the proof breaks down in the middle degrees $k=n / 2$, $(n \pm 1) / 2$.

As in the first step, $E_{1}$ is obtained by an iterative process. The first term $E_{1,0}$ is chosen to satisfy

$$
N_{p}\left(\Delta_{g}\right) N_{p}\left(E_{1,0}\right)=N_{p}\left(Q_{0}\right)
$$

for each $p \in \partial M$. In fact, by (3.47) there is a function

$$
N_{p}\left(E_{1,0}\right) \in A_{p h g}^{\sigma^{\prime}, \tau^{\prime}}\left(F_{p}\right) \otimes \Gamma_{0}^{1 / 2}
$$

solving this equation, and then, by the exactness of the sequence in (2.45) there is an $E_{1,0}$ with this function as its normal operator. Thus

$$
\Delta_{g^{\prime}} E_{1,0}-Q_{0}=-Q_{1,1}^{\prime} \in R \Psi_{0}^{-\infty, \sigma^{\prime}, \tau^{\prime}}
$$

It is crucial to observe now that $Q_{1,1}^{\prime}$ is actually slightly better:

$$
\mathrm{Q}_{1,1}^{\prime} \in \mathrm{R} \Psi_{0}^{-\infty}, \sigma^{\prime}+1, \tau^{\prime}
$$

$$
\Rightarrow Q_{1,1}^{\prime}=x Q_{1,1}^{\prime \prime}, \quad Q_{1,1}^{\prime \prime} \in \Psi_{0}^{-\infty, \sigma^{\prime}, \tau^{\prime}}
$$

The reason for this gain is that the normal operator of $\Delta_{g}$ annihilates the top order term of the asymptotic expansion for $E_{1,0}$ at the top face. Indeed, this is so because $N_{p}\left(E_{1,0}\right)$ is given by applying the Green function for this normal operator to $N_{p}\left(Q_{0}\right)$ (which is supported in the interior), and now (3.47) implies that $E_{1,0}$ has the correct top order term.

To get further terms in the series for $E_{1}$, let us assume that, instead of a Taylor series in $R$, we let

$$
E_{1} \sim \sum \tilde{\mathbf{x}}^{\mathbf{j}} \mathrm{E}_{1, \mathrm{j}}
$$

so that

$$
\Delta_{g} E_{1} \sim \sum \tilde{x}^{j_{j}}{ }_{\mathbf{g}} E_{1, j}
$$

Define $Q_{1, j}=\underset{\sim}{\underset{\sim}{x}} Q_{1, j}^{\prime \prime} ;$ then the inductive step is to solve

$$
\Delta_{\mathrm{g}} \mathrm{E}_{1, j}-Q_{1, j}=-x Q_{1, j+1}^{\prime \prime}
$$

where

$$
Q_{1, j} \in \Psi_{0}^{-\infty, \sigma^{\prime}+1, \tau^{\prime}-j} \Leftrightarrow Q_{1, j}^{\prime \prime} \in \Psi_{0}^{-\infty}, \sigma^{\prime}, \tau^{\prime-j+1} .
$$

This has already been accomplished for $j=0$. For greater values of $j$, we solve this equation by first solving

$$
N_{p}\left(\Delta_{g}\right) N_{p}\left(E_{1, j}\right)=N_{p}\left(Q_{1, j}\right) \in A^{\sigma^{\prime}+1, \tau^{\prime}-j}\left(F_{p} ;^{o} \Lambda^{k_{\otimes}} \otimes \Gamma_{0}^{1 / 2}\right) .
$$

By Theorem (3.50) an appropriate solution may indeed be found; its extension into the interior of $M \times{ }_{0} M$ then satisfies

$$
\begin{gathered}
\Delta_{g} E_{1, j}-Q_{1, j}=-R Q_{1, j+1}^{\prime} \\
Q_{1, j+1}^{\prime} \in \Psi_{0}^{-\infty}, \sigma^{\prime}, \tau^{\prime-j}
\end{gathered}
$$

However, by reasoning as above -and now the point to note is that by hypothesis the top order term in the expansion for $Q_{1, j}$ is of type $\sigma^{\prime}+1$ and so doesn't interfere with the argument -we actually have

$$
\begin{gathered}
\Delta_{g^{\prime}} E_{1, j}-Q_{1, j}=-x Q_{1, j+1}^{\prime \prime} \\
Q_{1, j+1}^{\prime \prime} \in \Psi_{0}^{-\infty, \sigma^{\prime}, \tau^{\prime}-j}
\end{gathered}
$$

as desired. Finally, since $E_{1, j} \in \Psi_{0}^{-\infty, \sigma^{\prime}, \tau^{\prime}-j}$

$$
\tilde{\mathbf{x}}^{\mathbf{j}^{E}}{ }_{1, \mathrm{j}} \in \Psi_{0}^{-\infty, \sigma^{\prime}, \tau^{\prime}}
$$

and

$$
E_{1} \sim \sum \tilde{x}^{j} E_{1, j} \in \Psi_{0}^{-\infty, \sigma^{\prime}, \tau^{\prime}}
$$

This $E_{1}$ solves (4.2).
The final step of this construction is to find $E_{2} \epsilon$ $\mathrm{R}^{\infty} \Psi^{-\infty}, \sigma^{\prime}, \tau^{\prime}$ such that

$$
\begin{equation*}
\Delta_{\mathrm{g}} \mathrm{E}_{2}-\mathrm{Q}_{1}=-\mathrm{Q}_{2} \in \mathrm{R}^{\infty} \Psi_{0}^{-\infty, \infty, \tau^{\prime}} \tag{4.3}
\end{equation*}
$$

This is quite easy: one merely solves for each term in the expansion at $T$ using the indicial operator, cf. the proof of (3.47). This completes the construction.

## Set

$$
E=E_{1}+E_{2}+E_{3}
$$

(4.4) Theorem: $E \in \Psi_{0}^{-2, \sigma^{\prime}, \tau^{\prime}}\left(M ;{ }^{o} \Lambda^{k} \otimes \Gamma_{0}^{1 / 2}\right)$ is a parametrix for $\quad \Delta_{g}$ :

$$
\Delta_{g} E=I-Q, Q \in R^{\infty} \Psi_{0}^{-\infty, \infty}, \tau^{\prime}\left(M ;^{o} \Lambda^{k} \otimes \Gamma_{0}^{1 / 2}\right)
$$

(4.5) Corollary: $\mathrm{Q}: \mathrm{L}^{2}\left({ }^{\mathrm{o}} \Lambda^{\mathrm{k}} \otimes \Gamma_{0}^{1 / 2}\right) \longrightarrow \mathrm{L}^{2}\left({ }^{\mathrm{o}} \Lambda^{\mathrm{k}} \otimes \Gamma_{0}^{1 / 2}\right)$ is compact, and thus

$$
\Delta_{\mathrm{g}}: \mathrm{L}^{2}\left({ }^{\mathrm{o}} \Lambda^{\mathrm{k}} \otimes \Gamma_{0}^{1 / 2}\right) \longrightarrow \mathrm{L}^{2}\left({ }^{\mathrm{o}} \Lambda^{\mathrm{k}} \otimes \Gamma_{0}^{1 / 2}\right) \text { is Fredholm. }
$$

## B. Proof of the Main Theorem

It follows from Corollary (4.5) that the harmonic space

$$
\mathscr{H}^{k}=\left\{\omega \in L^{2} \Omega^{k}(\mathrm{dg}): \Delta_{\mathrm{g}} \omega=0\right\}
$$

is finite dimensional. It remains for us to identify it in terms of the topological cohomology of $M$. First observe that the adjoint of the remainder

$$
Q^{*} \in R^{\infty} \Psi_{0}^{-\infty}, \sigma^{\prime},{ }^{\infty}\left(M ;{ }^{o}{ }^{k} \otimes \Gamma_{0}^{1 / 2}\right)
$$

To double check the order of $\kappa\left(Q^{*}\right)$ on the top face, notice that

$$
E^{*} \Delta_{g}=I-Q^{*}
$$

implies

$$
\begin{equation*}
\omega \in \forall^{k} \Rightarrow Q^{*}(\omega \cdot \mu)=\omega \cdot \mu \tag{4.6}
\end{equation*}
$$

Thus the order $\sigma^{\prime}$ is the only one possible which ensures that

$$
I\left(\Delta_{g}\right) \omega \equiv 0
$$

modulo terms of one order higher.
Remove the density factor and regard $\omega$ as a section of the ordinary form bundle. From (2.55) and (4.6)

$$
\omega \in \mathscr{H},\left|k-\frac{n}{2}\right|>\frac{1}{2} \Rightarrow \omega \in A_{\mathrm{phg}}^{\alpha}\left(M, \Lambda^{\mathrm{k}}\right)
$$

(4.7)

$$
\alpha=\left\{\begin{array}{cl}
(n-2 k-1, n-2 k) & k<\frac{n-1}{2} \\
(0,-1) & k>\frac{n+1}{2}
\end{array} .\right.
$$

These explicit asymptotics of $L^{2}$ harmonic forms, together with various consequences of the existence of the parametrix, allow us to identify $\mathscr{H}^{k}$ with the deRham cohomology spaces of $M$.

Let us first briefly review the deRham theory on a compact manifold with boundary. There are, of course, two flavours of cohomology to consider, the relative and the absolute, and these are dual. Inasmuch as $M$ and its interior share the same topological cohomology, it is reasonable to let the absolute cohomology be defined by

$$
\mathrm{H}^{\mathrm{k}}(\mathrm{M})=\left\{\omega \in \Omega^{\mathrm{k}}(\stackrel{\circ}{\mathrm{M}}): \mathrm{d} \omega=0\right\} /\left\{\mathrm{d} \eta: \eta \in \Omega^{\mathrm{k}-1}(\stackrel{\circ}{\mathrm{M}})\right\}
$$

In this equality we must specify what sort of regularityand boundary values-these forms should possess. In fact. it suffices either to let $\omega$ and $\eta$ be $C^{\infty}$ in ${ }^{M}$ or just distributional (currents), cf. [9], with no boundary constraints. However, when proving the Hodge theorem for a nondegenerate metric it is more convenient to use smooth forms on $M$ which satisfy 'absolute' boundary conditions, cf. [26], which we shall not define here since we do not need them. On the other hand, the relative cohomology not surprisingly always requires constraints at the boundary. Hence, it is well known that we may compute $H^{k}(M, \partial M)$ using the complex of forms satisfying relative boundary conditions

$$
H^{k}(M, \partial M)=\frac{\left\{\omega \in C^{\infty} \Omega^{k}(M): i^{*} \omega=0, d \omega=0\right\}}{\left\{d \eta: \eta \in C^{\infty} \Omega^{k-1}(M), i^{*} \omega=0\right\}}
$$

Here $i: \partial M \longrightarrow M$ is the inclusion. We could equally well use forms supported in the interior of $M$, as shown by an argument using the chain homotopy operator in the proof of our next theorem.

We need to define the relative groups using somewhat less regular forms. There are two spaces of distributions on $M$, briefly mentioned already in Chapter 2, but see [22] for more details. Let $\dot{C}^{\infty}$ denote the smooth functions vanishing to infinite order at $\partial M$ and let

$$
C^{-\infty}(M)=\left(\dot{C}^{\infty}(M)\right)^{\prime}, \quad \dot{C}^{-\infty}(M)=\left(C^{\infty}(M)\right)^{\prime}
$$

These are the 'extendible' and 'supported' distributions, respectively. If $\tilde{M}$ is an open extension of $M$, then the former space contains the restrictions to $\stackrel{\circ}{\mathrm{M}}$ of arbitrary distributions on $\tilde{M}$, while elements of the latter are distributions on $\tilde{M}$ actually supported in M. The sequence

$$
0 \longrightarrow \mathrm{C}^{-\infty}(M, \partial M) \longrightarrow \dot{C}^{-\infty}(M) \longrightarrow \mathrm{C}^{-\infty}(M) \longrightarrow 0
$$

is exact, the initial space being of distributions on $M$ supported in $\partial M$. One of the few differences between these spaces arises by comparing how differential operators act on them. As may be checked from duality, an extendible distribution is to be differentiated in $\stackrel{\circ}{M}$, the result extending to the boundary by continuity; a supported distribution is differentiated as a distribution on $\tilde{M}$, so that boundary layers may well occur.

Having stated these facts, we now come to the point of the discussion. It is the case (which unfortunately is not to be found in the literature) that the absolute and relative cohomologies of $M$ may be computed from the complexes of forms with values in $C^{-\infty}$ and $\dot{C}^{-\infty}$, respectively. We mention all this to justify a minor technicality concerning regularity in the proof of our
(4.8) Main Theorem: There are natural isomorphisms

$$
\begin{aligned}
\mathscr{H}^{\mathrm{k}} \ni \omega \longmapsto[\omega] \in \mathrm{H}^{\mathrm{k}}(M, \partial M) & \mathrm{k}<\frac{\mathrm{n}-1}{2} \\
\mathscr{H}^{\mathrm{k}} \ni \omega \longmapsto[\omega] \in \mathrm{H}^{\mathrm{k}}(M) & \mathrm{k}>\frac{\mathrm{n}+1}{2} .
\end{aligned}
$$

Proof: $L^{2}$ harmonic forms are closed and coclosed since the metric $g$ is complete. Hence when $k>\frac{\mathrm{n}+1}{2}, \omega \in \mathscr{H}^{\mathrm{k}}$ represents an absolute cohomology class. If $k<\frac{n-1}{2}$, then from (4.7) both components of $\omega$ vanish at $\partial M$. We may then consider it as lying in $\dot{\mathrm{C}}^{-\infty}(M)$, and it is still closed in the sense of supported distributions. Thus $\omega$ represents a relative class. It will suffice to show that below the middle degree the map is infective, while above the middle degree it is surjective.

Suppose then that for $k<\frac{n-1}{2}, \omega \in \mathscr{H}^{k}$, we have

$$
\omega=\mathrm{d} \eta, \quad \mathrm{i}^{*} \eta=0 .
$$

Choose finitely many coordinate patches near $\partial M$ such that $h_{\alpha n}=0, h_{n n}=1$, and the coordinate $z^{n}$ is defined globally near the boundary. Then, since $\delta_{g} \omega=0$, we compute

$$
\begin{aligned}
\|\omega\|^{2} & =\langle\mathrm{d} \eta, \omega\rangle-\left\langle\eta, \delta_{\mathrm{g}} \omega\right\rangle=\int_{\mathrm{M}} \mathrm{~d}\left(\eta \wedge{ }_{\mathrm{g}} \omega\right) \\
& =\lim _{\epsilon \rightarrow 0} \int_{\left\{\mathrm{z}^{\mathrm{n}}=\epsilon\right\}} \eta \wedge *_{\mathrm{h}} \omega \rho^{2 \mathrm{k}-\mathrm{n}} .
\end{aligned}
$$

By assumption on the coordinates, only $\eta_{t}$ and $\omega_{n}=$ $\left({ }^{*}{ }^{\omega}{ }^{\omega}\right)_{t}$ enter into this last integral. We know that

$$
\omega_{\mathrm{n}}=O\left(\rho^{\mathrm{n}-2 k} \log \rho\right)
$$

but it is less clear at what rate $\eta_{t}$ vanishes.
We study this issue by using the chain homotopy operator of [26]. Thus if

$$
\omega=\omega_{I} d y^{I}+\omega_{J} d y^{J} \wedge d x \in \Omega^{k}
$$

and $\omega_{I}(y, 0)=0$, then we define $R \omega=(R \omega)_{J} d y$ by

$$
(R \omega)_{J}=(-1)^{k-1} \int_{0}^{1} \omega_{J}(y, t x) x d t
$$

This is actually coordinate independent, and satisfies

$$
\mathrm{dR} \omega+\mathrm{Rd} \omega=\omega
$$

so long as $i^{*} \omega=0$. Note also that $i^{*} R \omega=0$.

If $\omega$ is closed then $\omega=\mathrm{dR} \omega$ in the neighbourhood of the boundary where $R \omega$ is defined. If $\psi$ is a function supported in this neighbourhood which is identically one near $\partial M$, then

$$
\omega-\mathrm{d}(\psi \mathrm{R} \omega)
$$

represents the same relative cohomology class and vanishes near the boundary. We are assuming that this class is trivial, so there exists a form $\beta$ for which

$$
\mathrm{i}^{*} \beta=0, \quad \omega-\mathrm{d}(\psi \mathrm{R} \omega)=\mathrm{d} \beta
$$

which is $C^{\infty}$ up to $\partial M$. This shows that we may take

$$
\eta=\beta+\psi \mathrm{R} \omega
$$

and so

$$
\eta_{t}=O(\rho \log \rho) .
$$

The limit above must now vanish, hence $\omega=0$ and the map into relative cohomology is injective.

On the other hand, when $k>\frac{n+1}{2}$, take a smooth representative $\alpha$ for an arbitrary absolute cohomology class. Obviously $\alpha$ is square-integrable with respect to
$\rho^{2 k-n}$ since $|\alpha|_{h}^{2}$ is bounded. We know that $\Delta_{g}$ has closed range with a finite dimensional complement, so we may write

$$
\alpha=\Delta_{g} \beta+\omega
$$

for some $\beta \in L^{2} \Omega^{k}$ and $\omega \in \mathscr{H}^{k}$. Obviously $\beta$ is $C^{\infty}$ in the interior of $M$, and since $d \alpha=0$ it follows that

$$
\alpha=\mathrm{d} \delta_{g} \beta+\omega
$$

Hence $[\alpha]=[\omega]$ and the theorem is proved.

As the last part of this proof shows, we have actually proved something about the $L^{2}$ cohomology spaces discussed in the introduction.
(4.9) Proposition: For $k<\frac{n-1}{2}$ or $k>\frac{n+1}{2}$ we have

$$
L^{2} H^{k} \simeq \mathscr{H}^{k}
$$

and so, in these degrees, $L^{2} H^{k}$ is finite dimensional.

Proof: We have established a strong Hodge decomposition for these values of $k$ :

$$
L^{2} \Omega^{k}(M, d g) \quad \alpha=\mathrm{d} \delta_{g} \beta+\delta_{g} \mathrm{~d} \beta+\omega
$$

$\omega \in \mathscr{H}^{k} . \alpha$ is closed iff the second summand vanishes, and so $[\alpha]=[\omega]$ in $L^{2} H^{k}$.

The existence of a parametrix away from the middle degrees implies that the spectrum of $\Delta_{g}$ on $L^{2} \Omega^{k}$ is bounded below by a positive constant for these values of k. However, when $k=(n \pm 1) / 2$-so that the dimension $n$ is odd-then judging from the special case $\mathbb{H}^{\mathbf{n}}$, cf. [11], the continuous spectrum should in general extend down to zero, and $\Delta_{g}$ would not have closed range. About the final case, when $n$ is even and $k=n / 2$, we can reach a definite conclusion.
(4.10) Proposition: $\mathscr{H}^{n / 2}$ is of infinite dimension, but on the orthocomplement $\left(\mathcal{H}^{\mathrm{n} / 2}\right)^{\perp}$ the spectrum of $\Delta_{g}$ has a positive lower bound. Hence the range of the Laplacian on $L^{2} \Omega^{n / 2}$ is closed, a strong Hodge decomposition holds, and $L^{2} H^{n / 2} \simeq \mathscr{H}^{\mathrm{n} / 2}$.

Proof: Since $L^{2} \Omega^{n / 2}(M, d g)=L^{2} \Omega^{n / 2}(M, d h) \quad$ and $\quad \delta_{g}=\rho^{2} \delta_{h}$ on $\Omega^{\mathrm{n} / 2}$ (the space $\mathscr{H}^{\mathrm{n} / 2}$ depends only on the conformal class of the metric), it is necessary only to solve the elliptic boundary problem

$$
\Delta_{h} \omega=0, \quad i^{*} \omega=\mu \in C^{\infty} \Omega^{n / 2}(\partial M), \quad i^{*} \delta_{h} \omega=0
$$

for the nondegenerate metric $h$; each solution is $C^{\infty}$, hence in $L^{2}(d g)$, and an infinite dimensional subspace of $H^{\mathrm{n} / 2}$ is thereby obtained.

The argument for the next part follows [14] closely. Choose $\epsilon>0$ so that $[0, \epsilon]$ does not intersect the spectrum of $\Delta_{g}^{n / 2 \pm 1}$ (the superscript indicating the degree of form on which $\Delta_{g}$ is supposed to act). Now suppose

$$
\omega=x\left(\Delta_{\mathrm{g}}^{\mathrm{n} / 2}\right) \alpha \neq 0
$$

By the assumption on $\epsilon$,

$$
\begin{gathered}
\mathrm{d} \omega=\mathrm{d} x\left(\Delta_{\mathrm{g}}^{\mathrm{n} / 2}\right) \alpha=x\left(\Delta_{\mathrm{g}}^{\mathrm{n} / 2+1}\right) \mathrm{d} \alpha=0 \\
\delta_{\mathrm{g}} \omega=\delta_{\mathrm{g}} \times\left(\Delta_{\mathrm{g}}^{\mathrm{n} / 2}\right) \alpha=x\left(\Delta_{\mathrm{g}}^{\mathrm{n} / 2-1}\right) \delta_{\mathrm{g}} \alpha=0
\end{gathered}
$$

and so $\omega \in \mathscr{H}^{\mathrm{n} / 2}$. This proves that

$$
\operatorname{spec}\left(\Delta_{g}^{n / 2}\right) \cap[0, \epsilon]=\{0\}
$$

as desired.
The last assertions are now immediate, and we are done.

It seems quite likely that a fairly complete spectral picture of the Hodge Laplacian for a conformally compact metric in any degree could be obtained from refinements of the analysis in this dissertation. Perhaps this issue will be addressed later if it seems warranted.

1. Anderson, M.: The Dirichlet Problem at Infinity for Manifolds of Negative Curvature, J. Diff. Geom., 18, no. 4 (1983), 701-721.
2. Anderson, M.: $L^{2}$ Harmonic Forms and a Conjecture of Dodziuk-Singer, Bulletin of the A.M.S. 13, no. 2 (1985), 163-165.
3. Anderson, M. and R. Schoen: Positive Harmonic Functions on Complete Manifolds of Negative Curvature; Preprint.
4. Andreotti, A. and E. Vesentini: Carleman Estimates for the Laplace-Beltrami Equation on Complex Manifolds, Publications Mathématiques IHES, no. 25 (1965).
5. Atiyah, M. F.: Elliptic Operators, Discrete Groups and von Neumann Algebras, Asterisque, 32-33 (1976), 43-72.
6. Bourguignon, J.-P. and H. Karcher: Curvature Operators: Pinching Estimates and Geometric Examples, Ann. Scient. Ec. Norm. Sup., 11 1978), 71-92.
7. Choi, H. I.: Thesis, University of California, Berkeley, 1982.
8. Coddington, E. and N. Levinson: Theory of Ordinary Differential Equations, McGraw Hill, New York (1955).
9. deRham, G.: Differentiable Manifolds, SpringerVerlag (1984).
10. Dodziuk, J.: $L^{2}$ Harmonic Forms on Rotationally Symmetric Riemannian Manifolds, Proc. Am. Math. Soc. 77 (1979), 395-400.
11. Donnelly, H.: The Differential Form Spectrum of Hyperbolic Space, Manuscripta Math. 33, (1981), 365-385.
12. Donnelly, H. and C. Fefferman: $L^{2}$-Cohomology and Index Theorem for the Bergman Metric, Ann. of Math., 118 (1983), 593-618.
13. Donnelly, H. and P. Li: Pure Point Spectrum and Negative Curvature for Noncompact Manifolds, Duke Math. J., 46, no. 3 (1979), 497-503.
14. Donnelly, H. and F. Xavier: On the Differential Form Spectrum of Negatively Curved Riemannian Manifolds, Amer. J. Math., 106 (1984), 169-185.
15. Eberlein, P. and B. O'Neill: Visibility Manifolds, Pacific J. Math., 46, no. 1 (1973), 45-109.
16. Fefferman, C.: The Bergman Kernel and Biholomorphic Mappings of Pseudoconvex Domains, Invent. Math. 26 (1974), 1-65.
17. Geroch, R.: Positive Sectional Curvature does not Imply Positive Gauss-Bonnet Integrand, Proc. A.M.S. Vol. 54 (1976), 267-270.
18. Hörmander, L.: The Analysis of Linear Partial Differential Operators, Vol. III, SpringerVerlag (1985).
19. Lebedev, N. N.: Special Functions and their Applications, Dover (1972).
20. Mazzeo, R. and R. B. Melrose: Meromorphic Extension of the Resolvent on Complete Spaces with Asymptotically Constant Negative Curvature, Preprint.
21. McKean, H.P.: An Upper Bound to the Spectrum of $\Delta$ on a Manifold of Negative Curvature, J. Diff. Geom., 4 (1970), 359-366.
22. Melrose, R. B.: Transformation of Boundary Problems Acta Math., 147 (1981), 149-236.
23. Melrose, R. B. and G. Mendoza: Elliptic Operators of Totally Characteristic Type, to Appear in J. Diff. Eq.
24. Melrose, R. B. and G. Mendoza: Degenerate Elliptic Boundary Problems, In Preparation.
25. Melrose, R. B. and N. Ritter: Interaction of Progressing Waves for Semilinear Wave Equations II, to Appear in Ark. Math.
26. Ray, D. and I. Singer: R-Torsion and the Laplacian on Riemannian Manifolds, Adv. in Math., 7 (1971). 145-210.
27. Sullivan, D.: The Dirichlet Problem at Infinity for a Negatively Curved Manifold, J. Diff. Geom., 18, no. 4 (1983), 723-732.
28. Yau, S. T.: Harmonic Functions on Complete Riemannian Manifolds, Comm. Pure Appl. Math., 28 (1975), 201-228.
29. Adams, R.A.: Sobolev Spaces, Academic Press, New York (1975).
