# Lectures on <br> Partial Differential Equations 

## By

G.B. Folland

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# Lectures on <br> Partial Differential Equations 

By<br>G.B. Folland<br>Lectures delivered at the<br>Indian Institute of Science, Bangalore

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## Preface

This book consists of the notes for a course I gave at the T.I.F.R. Center in Bangalore from September 20 to November 20, 1981. The purpose of the course was to introduce the students in the Programme in Application of Mathematics to the applications of Fourier analysis-by which I mean the study of convolution operators as well as the Fourier transform itself-to partial differential equations. Faced with the problem of covering a reasonably broad spectrum of material in such a short time, I had to be selective in the choice of topics. I could not develop any one subject in a really thorough manner; rather, my aim was to present the essential features of some techniques that are well worth knowing and to derive a few interesting results which are illustrative of these techniques. This does not mean that I have dealt only with general machinery; indeed, the emphasis in Chapter 2 is on very concrete calculation with distributions and Fourier transforms-because the methods of performing such calculations are also well worth knowing.

If these notes suffer from the defect of incompleteness, they posses the corresponding virtue of brevity. They may therefore be of value to the reader who wishes to be introduced to some useful ideas without having to plough through a systematic treatise. More detailed accounts of the subjects discussed here can be found in the books of Folland [1], Stein [2], Taylor [3], and Treves [4].

No specific knowledge of partial differential equations or Fourier Analysis is presupposed in these notes, although some prior acquittance with the former is desirable. The main prerequisite is a familiarity with the subjects usually gathered under the rubic "real analysis": measure
and integration, and the elements of point set topology and functional analysis. In addition, the reader is expected to be acquainted with the basic facts about distributions as presented, for example, in Rudin [7].

I wish to express my gratitude to professor K.G. Ramanathan for inviting me to Bangalore, and to professor S.Raghavan and the staff of the T.I.F.R. Center for making my visit there a most enjoyable one. I also wish to thank Mr S. Thangavelu and Mr K.T. Joseph for their painstaking job of writing up the notes.

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## Chapter 1

## Preliminaries

IN THIS CHAPTER, we will study some basic results about convolu- 1 tion and the Fourier transform.

## 1 General Theorems About Convolutions

We will begin with a theorem about integral operators.
Theorem 1.1. Let $K$ be a measurable function on $\mathbb{R}^{n} \times \mathbb{R}^{n}$ such that, for some $c>0$,

$$
\int|K(x, y)| d y \leq c, \int|K(x, y)| d x \leq c, \text { for every } x, y \text { in } \mathbb{R}^{n} .
$$

If $1 \leq p \leq \infty$ and $f \in L^{p}\left(\mathbb{R}^{n}\right)$, then the function $T f$, defined by $T f(x)=\int K(x, y) f(y) d y$ for almost every $x$ in $\mathbb{R}^{n}$, belongs to $L^{p}\left(\mathbb{R}^{n}\right)$ and further,

$$
\|T f\|_{p} \leq c\|f\|_{p}
$$

Proof. If $p=\infty$, the hypothesis $\int|K(x, y)| d x \leq c$ is superfluous and the conclusion of the theorem is obvious. If $p<\infty$, let $q$ denote the conjugate exponent. Then, by Hölder 's inequality,

$$
|T f(x)| \leq\left\{\int|K(x, y)| d y\right\}^{1 / q}\left\{\int|K(x, y)||f(y)|^{p} d y\right\}^{1 / p}
$$

$$
\leq c^{1 / q}\left\{\int|K(x, y)||f(y)|^{p} d y\right\}^{1 / p}
$$

From this we have,

$$
\begin{aligned}
\int|T f(x)|^{p} d x & \leq c^{p / q} \iint|K(x, y) \| f(y)|^{p} d y d x \\
& \leq c^{1+p / q} \int|f(y)|^{p} d x=c^{1+p / q}\|f\|_{p}^{p}
\end{aligned}
$$

$2 \quad$ Therefore $\|T f\|_{p} \leq c\|f\|_{p}$.
Next, we define the convolution of two locally integrable functions.
Definition 1.2. Let $f$ and $g$ be two locally integrable functions. The convolution of $f$ and $g$, denoted by $f * g$, is define by

$$
(f * g)(x)=\int f(x-y) g(y) d y=\int f(y) g(x-y) d y=(g * f)(x)
$$

provided that the integrals in question exist.
The basic theorem on convolution is the following theorem, called Young's inequality.

Theorem 1.3 (Young's Inequality). Let $f \in L^{1}\left(\mathbb{R}^{n}\right)$ and $g \in L^{p}\left(\mathbb{R}^{n}\right)$, for $1 \leq p \leq \infty$. Then $f * g \in L^{p}\left(\mathbb{R}^{n}\right)$ and

$$
\|f * g\|_{p} \leq\|g\|_{p}\|f\|_{1} .
$$

Proof. Take $K(x, y)=f(x-y)$ in Theorem 1.1. Then $K(x, y)$ satisfies all the conditions of Theorem 1.1 and the conclusion follows immediately.

The next theorem underlies one of the most important uses of convolution. Before coming to the theorem, let us prove the following

Lemma 1.4. For a function $f$ defined on $\mathbb{R}^{n}$ and $x$ in $\mathbb{R}^{n}$, we define a function $f^{x}$ by $f^{x}(y)=f(y-x)$. If $f \in L^{p}, 1 \leq p<\infty$, then $\lim _{x \rightarrow O} \| f^{x}-$ $f \|_{p}=0$.

3 Proof. If $g$ is a compactly supported continuous function, then $g$ is uniformly continuous, and so $g^{x}$ converges to $g$ uniformly as $x$ tends to 0 .

Further, for $|x| \leq 1, g^{x}$ and $g$ are supported in a common compact set. Therefore, $\lim _{x \rightarrow 0}\left\|g^{x}-g\right\|_{p}=0$. Given $f \in L^{p}$, we can find a function $g$ which is continuous and compactly supported such that $\|f-g\|_{p}<\in / 3$ for $\in>0$. But then $\left\|g^{x}-f^{x}\right\|_{p}<\in / 3$ also holds. Therefore

$$
\begin{aligned}
\left\|f^{x}-f\right\|_{p} & \leq\left\|f^{x}-g^{x}\right\|_{p}+\left\|g^{x}-g\right\|_{p}+\|g-f\|_{p} \\
& \leq 2 \in / 3+\left\|g^{x}-g\right\|_{p}
\end{aligned}
$$

Since $\lim x \rightarrow 0\left\|g^{x}-g\right\|_{p}=0$, we can choose $x$ close to 0 so that $\left\|g^{x}-g\right\|_{p}<\epsilon / 3$. Then $\left\|f^{x}-f\right\|<\epsilon$ and this proves the lemme since $\in$ is arbitrary.

Remark 1.5. The above lemma is false when $p=\infty$. Indeed, " $f^{x} \rightarrow f$ in $L^{\infty "}$ means precisely that $f$ is uniformly continuous.

Let us now make two important observations about convolutions which we shall use without comment later on.
i) $\operatorname{Supp}(f * g) \subset \operatorname{Supp} f+\operatorname{Supp} g$, where

$$
A+B=\{x+1: x \in A, y \in B\} .
$$

ii) If $f$ is of class $C^{k}$ and $\partial^{\alpha}(|\alpha| \leq k)$ and $g$ satisfy appropriate conditions so that differentiation under the integral sign is justified, then $f * g$ is of class $C^{k}$ and $\partial^{\alpha}(f * g)=\left(\partial^{\alpha} f\right) * g$.

Theorem 1.6. Let $g \in L^{1}\left(\mathbb{R}^{n}\right)$ and $\int g(x) d x=$ a. Let $g \in(x)=\epsilon^{-n} g(x / \in) \quad 4$ for $\in>0$. Then, we have the following:
i) If $f \in L^{p}\left(\mathbb{R}^{n}\right), p<\infty, f * g_{\in}$ converges to af in $L^{p}$ as tends to 0 .
ii) If $f$ is bounded and continuous, then $f * g_{\in}$ converge to af uniformly on compact sets as tends to 0 .

Proof. By the change of variable $x \rightarrow \in x$, we see that $\int g_{\epsilon}(x) d x=a$ for all $\in>0$. Now

$$
\begin{aligned}
\left(f * g_{\epsilon}\right)(x)-a f(x) & =\int f(x-y) g_{\epsilon}(y) d y-\int f(x) g_{\epsilon}(y) d y \\
& =\int[f(x-y)-f(x)] g_{\in}(y) d y \\
& =\int[f(x-\in y)-f(x)] g(y) d y \\
& =\int\left[f^{\in y}(x)-f(x)\right] g(y) d y
\end{aligned}
$$

If $f \in L^{p}$ and $p<\infty$, we apply Minkoswski's inequality for integrals to obtain

$$
\left\|f * g_{\epsilon}-a f\right\|_{p} \leq \int\left\|f^{\in y}-f\right\|_{p}|g(y)| d y
$$

The function $y \rightarrow\left\|f^{\in y}-f\right\|_{p}$ is bounded by $2\|f\|_{p}$ and tends to 0 as $\in$ tends to 0 for each $y$, by lemma 1.4. Therefore, we can apply Lebesgue Dominated Convergence theorem to get the desired result.

On the other hand, suppose $f$ is bounded and continuous. Let $K$ be any compact subset of $\mathbb{R}^{n}$. Given $\delta>0$, choose a compact set $G \mathbb{R}^{n}$ such that

$$
\int_{\mathbb{R}^{n}-G}|g(y)| d y<\delta .
$$

Then

$$
\begin{aligned}
\operatorname{Sup}_{x \in k}\left(\mid f * g_{\epsilon}\right)(x)-a f(x) \mid & \leq 2 \delta\|f\|_{\infty} \\
& +\operatorname{Sup}_{(x, y) \epsilon K \times G}|f(x-\epsilon y)-f(x)| \int_{G}|g| d y .
\end{aligned}
$$

Since $f$ is uniformly continuous on the compact set $K$ the second term tends to 0 as $\in$ to 0 . Since $\delta$ is arbitrary, we see that

$$
\sup x \in K\left|\left(f * g_{\epsilon}\right)(x)-a f(x)\right| \rightarrow 0 \text { as } \epsilon \rightarrow 0
$$

Hence the theorem is proved.

Corollary 1.7. The space $C_{o}^{\infty}\left(\mathbb{R}^{n}\right)$ is dense in $L^{p}\left(\mathbb{R}^{\propto}\right)$ for $1 \leq p<\infty$.
Proof. Let

$$
\begin{array}{rlrl}
\phi(x) & =e^{-1 /\left(1-|x|^{2}\right)} & \text { for }|x|<1 \\
& =0 & & \text { for }|x| \leq 1
\end{array}
$$

Then $\phi \epsilon C_{o}^{\infty}\left(\mathbb{R}^{n}\right)$ and $\int \phi(x) d x=1 / a>0$. If $f \epsilon L^{p}$ and has compact support, then $a\left(f * \phi_{\epsilon}\right) \epsilon C_{o}^{\infty}\left(\mathbb{R}^{n}\right)$ and by theorem 1.6, $a\left(f * \phi_{\epsilon}\right)$ converges to $f$ in $L^{p}$ as $\epsilon$ tends to 0 . Since $L^{p}$ functions with compact support are dense in $L^{p}$, this completes the proof.

Proposition 1.8. Suppose $K \subset \mathbb{R}^{n}$ is compact and $\Omega \supset K$ be an open subset of $\mathbb{R}^{n}$. Then there exists a $C_{o}^{\infty}$ function $\phi$ such that $\phi(x)=1$ for $x \in K$ and $\operatorname{Supp} \phi \subset \Omega$.

Proof. Let $V=\left\{x \in \Omega: d(x, K) \leq \frac{1}{2} \delta\right\}$ where $\delta=d\left(K, \mathbb{R}^{n} \backslash \Omega\right)$. Choose 6 a $\phi_{o} \epsilon C_{o}^{\infty}$ such that $\operatorname{Supp} \phi_{o} \subset B\left(0, \frac{1}{2} \delta\right)$ and $\int \phi_{o}(x) d x=1$. Define

$$
\phi(x)=\int_{V} \phi_{o}(x-y) d y=\left(\phi_{0} * X_{V}\right)(x)
$$

Then $\phi(x)$ is a function with the required properties.

## 2 The Fourier Transform

In this section, we will give a rapid introduction to the theory of the Fourier transform.

For a function $f \in L^{1}\left(\mathbb{R}^{n}\right)$, the Fourier transform of the function $f$, denoted by $\hat{f}$, is defined by

$$
\hat{f}(\xi)=\int e^{-2 \pi i x \cdot \xi} f(x) d x, \xi \in \mathbb{R}^{n}
$$

Remark 1.9. Our definition of $\hat{f}$ differs from some other in the placement of the factor of $2 \pi$.

BASIC PROPERTIES OF THE FOURIER TRANSFORM
For

$$
\begin{equation*}
f \epsilon L^{1},\|\hat{f}\|_{\infty} \leq\|f\|_{1} . \tag{1.10}
\end{equation*}
$$

The proof of this is trivial.
For

$$
\begin{equation*}
f, g \in L^{1},(f * g)^{\wedge}(\xi)=\hat{f}(\xi) \hat{g}(\xi) \tag{1.11}
\end{equation*}
$$

Indeed,

$$
\begin{aligned}
(f * g)(\xi) & =\iint e^{-2 \pi i x \cdot \xi} f(y) g(x-y) d y d x \\
& =\iint e^{-2 \pi i(x-y) \cdot \xi} g(x-y) e^{-2 \pi i Y \cdot \xi} f(y) d y d x \\
& =\int e^{-2 \pi i(x-y) \cdot \xi} g(x-y) d x \int f(y) e^{-2 \pi i y \cdot \xi} d y \\
& =\hat{f}(\xi) \hat{g}(\xi)
\end{aligned}
$$

7
Let us now consider the Fourier transform in the Schwartz class $S=$ $S\left(\mathbb{R}^{n}\right)$.

Proposition 1.12. For $f \in S$, we have the following:
i) $\hat{f} \epsilon C^{\infty}\left(\mathbb{R}^{n}\right)$ and $\partial^{\beta} \hat{f}=\hat{g}$ where $g(x)=(-2 \pi i x)^{\beta} f(x)$.
ii) $\left(\partial^{\beta} f \hat{)}(\xi)=(2 \pi i \xi)^{\beta} \hat{f}(\xi)\right.$.

Proof. i) Differentiation under the integral sign proves this.
ii) For this, we use integration by parts.

$$
\begin{aligned}
\left(\partial^{\beta} f \hat{)(\xi)}\right. & =\int e^{-2 \pi i x \cdot \xi}\left(\partial^{\beta} f\right)(x) d x \\
& =(-1)^{|\beta|} \int \partial_{\beta}\left[e^{-2 \pi i x \cdot \xi}\right] f(x) d x \\
& =(-1)^{|\beta|}(-2 \pi i \xi)^{\beta} \int e^{-2 \pi i x \cdot \xi} f(x) d x \\
& =(2 \pi i \xi)^{\beta} \hat{f}(\xi) .
\end{aligned}
$$

Corollary 1.13. If $f \in S$, then $\hat{f} \in S$ also.
Proof. For multi-indices $\alpha$ and $\beta$, using proposition 1.12 we have

$$
\begin{aligned}
\xi^{\alpha}\left(\partial^{\beta} \hat{f}\right)(\xi) & =\xi^{\alpha}\left((-2 \pi i x)^{\beta} f(x)\right)^{\wedge}(\xi) \\
& =(2 \pi i)^{-|\alpha|}\left[\partial^{\alpha}\left((-2 \pi i x)^{\beta} f(x) \hat{)}(\xi)\right]\right. \\
& =(-1)^{|\beta|}(2 \pi i)^{|\beta|-|\alpha|}\left(\partial^{\alpha}\left(x^{\beta} f(x)\right)\right)^{\hat{}}(\xi)
\end{aligned}
$$

Since $f \in S, \partial^{\alpha}\left(x^{\beta} f(x)\right) \epsilon L^{1}$ and hence $\left(\partial^{\alpha}\left(x^{\beta} f(x)\right)\right)^{\wedge} \epsilon L^{\infty}$. Thus $\xi^{\alpha}$ $\left(\partial^{\beta} \hat{f}\right)$ is bounded. Since $\alpha$ and $\beta$ are arbitrary, this proves that $\hat{f} \epsilon S$.

Corollary 1.14 (RIEMANN-LEBESGUE LEMMA). If $f \epsilon L^{1}$, then $\hat{f}$ is continuous and vanishes at $\infty$.

Proof. By Corollary 1.13 this is true for $f \epsilon S$. Since $S$ is dense in $L^{1} 8$ and $\|\hat{f}\|_{\infty} \leq\|f\|_{1}$, the same in true for all $f \in L^{1}$.

Let us now compute the Fourier transform of the Gaussian.
Theorem 1.15. Let $f(x)=e^{-\pi a|x|^{2}}$, Re $a>0$. Then, $\hat{f}(\xi)=a^{-n / 2}$ $e^{-a^{-1}} \pi|\xi|^{2}$.

Proof.

$$
\hat{f}(\xi)=\int e^{-2 \pi i x \cdot \xi} e^{-a \pi|x|^{2}} d x
$$

i.e.,

$$
\hat{f}(\xi)=\prod_{j=1}^{n} \int_{-\infty}^{\infty} e^{-2 \pi i x_{j} \xi_{j}} e^{-a \pi x_{j}^{2}} d x_{j} .
$$

Thus it suffices to consider the case $n=1$. Further, we can take $a=1$ by making the change of variable $x \rightarrow a^{-1 / 2} x$.

Thus we are assuming $f(x)=e^{-\pi x^{2}}, x \in \mathbb{R}$. Observe that $f^{\prime}(x)+$ $2 \pi x f(x)=0$. Taking the Fourier transform, we obtain

$$
2 \pi i \xi \hat{f}(\xi)+i \hat{f}^{\prime}(\xi)=0
$$

Hence

$$
{\hat{f^{\prime}}}^{\prime}(\xi) / f(\xi)=-2 \pi \xi
$$

which, on integration, gives $f(\xi)=c e^{-\pi \xi^{2}}, c$ being a constant.
The constant $c$ is given by

$$
c=\hat{f}(0)=\int_{-\infty}^{\infty} e^{-\pi x^{2}} d x=1
$$

Therefore $\hat{f}(\xi)=e^{-\pi \xi^{2}}$, which completes the proof.
We now derive the Fourier inversion formula for the Schwartz class $S$.

Let us define $f^{V}(\xi)=\int e^{2 \pi i x \cdot \xi} f(x) d x=\hat{f}(-\xi)$.
Theorem 1.16 (Fouries Inversion Theorem). For $f \epsilon(\hat{f})^{\vee}=f$.
Proof. First, observe that for $f, g \epsilon L^{1}, \int f \hat{g}=\int \hat{f} g$. In fact,

$$
\begin{aligned}
\int \hat{f}(x) g(x) d x & =\iint e^{-2 \pi i y \cdot x} f(y) g(x) d y d x \\
& =\int\left[\int e^{-2 \pi i y \cdot x} g(x) d x\right] f(y) d y \\
& =\int \hat{g}(y) f(y) d y
\end{aligned}
$$

Given $\epsilon>0$ and $x$ in $\mathbb{R}^{n}$, take the function $\phi$ defined by

$$
\phi(\xi)=e^{-2 \pi i x \cdot \xi-\pi \epsilon^{2}|\xi|^{2}} .
$$

Now

$$
\begin{aligned}
\hat{\phi}(y) & =\int e^{-2 \pi i y \cdot \xi} e^{-2 \pi i x \cdot \xi-\pi \epsilon^{2} \mid \xi \xi^{2}} d \xi \\
& =\int e^{-2 \pi i(y-x) \cdot \xi \xi_{e}-\pi \epsilon^{2}|\xi|^{2}} d \xi \\
& =\epsilon^{-n} e^{-\pi \epsilon^{-2}|x-y|^{2}} .
\end{aligned}
$$

If we take $g(x)=e^{-\pi|x|^{2}}$ and define $g_{\epsilon}(x)=\epsilon^{-n} g(x / \epsilon)$, then

$$
\hat{\phi}(y)=g_{\epsilon}(x-y) .
$$

Therefore

$$
\begin{aligned}
\int e^{2 \pi i x \cdot \xi} \hat{f}(\xi) e^{-\pi \epsilon^{2}|\xi|^{2}} d \xi & =\int \hat{f}(\epsilon) \phi(\xi) d \xi \\
& =\int f(y) \hat{\phi}(y) d y \\
& =\int f(y) g_{\epsilon}(x-y) d y \\
& =\left(f * g_{\epsilon}\right)(x)
\end{aligned}
$$

But as $\epsilon$ tends to $0,\left(f * g_{\epsilon}\right)$ converges to $f$, by Theorem 1.6 and $\mathbf{1 0}$ clearly

$$
\int e^{2 \pi i \cdot x \xi} \hat{f}(\xi) e^{-\pi \epsilon^{2}|\xi|^{2}} d \xi \rightarrow \int \hat{f}(\xi) e^{2 \pi i \cdot x \cdot \xi_{d \xi}}
$$

Therefore $(\hat{f})^{v}=f$.
Corollary 1.17. The Fourier transform is an isomorphism of $S$ onto $S$.
Next, we prove the Plancherel Theorem.
Theorem 1.18. The Fourier transform uniquely extends to a unitary map of $L^{2}\left(\mathbb{R}^{n}\right)$ onto itself.

Proof. For $f \epsilon S$, define $\tilde{f}(x)=\overline{f(-x)}$. Then it is easily checked that $\hat{\tilde{f}}=\overline{\hat{f}}$, so that

$$
\begin{aligned}
\|f\|_{2}^{2} & =\int|f(x)|^{2} d x \\
& =\int f(x) \tilde{f}(-x) d x \\
& =(f * \tilde{f})(0) \\
& =\int\left(f_{*} \tilde{f}\right) \hat{)}(\xi) d \xi \\
& =\int \hat{f}(\xi) \hat{\tilde{f}}(\xi) d \xi \\
& =\int \hat{f}(\xi) \overline{\hat{f}}(\xi) d \xi=\|\hat{f}\|_{2}^{2}
\end{aligned}
$$

Therefore, the Fourier transform extends continuously to an isometry of $L^{2}$. It is a unitary transformation, since its image $S$ is dense in $L^{2}$ 。

Let us observe how the Fourier transform interacts with translations, rotations and dilations.
(1.19) The Fourier transform and translation: If $f^{x}(y) f(y-x)$ then

$$
\begin{aligned}
\hat{f}^{x}(\xi) & =\int e^{2 \pi i y \cdot \xi} f(y-x) d y \\
& =\int e^{2 \pi i(z+x) \cdot \xi} f(z) d z(\text { by putting } y-x=z) \\
& =\int e^{2 \pi i x \cdot \xi} \hat{f}(\xi)
\end{aligned}
$$

(1.20) The Fourier transform and rotations (orthogonal transformations):

Let $T: \mathbb{R}^{n} \rightarrow \mathbb{R}^{n}$ be an orthogonal transformation. Then

$$
\begin{aligned}
(f \circ T) \hat{(\xi)} & =\int e^{-2 \pi i x \cdot \xi}(f o T)(x) d x \\
& =\int e^{-2 \pi i T^{-1} y \cdot \xi} f(y) d y(\text { by putting } y=T x) \\
& =\int e^{-2 \pi i y \cdot T \xi} f(y) d y \\
& =\hat{f}(T \xi)=(\hat{f} o T)(\xi) .
\end{aligned}
$$

Thus, $(f o T)^{\wedge}=\hat{f o} T$ i.e. ${ }^{\wedge}$ commutes with rotations.
(1.21) The Fourier transform and dilation: Let $f_{r}(x)=r^{-n} f(x / r)$.

Then

$$
\begin{aligned}
\hat{f}_{r}(\xi) & =\int e^{-2 \pi i x \cdot \xi} r^{-n} f(x / r) d x \\
& =\int e^{-2 \pi i y \cdot \xi} f(y) d y=\hat{f}(r \xi)
\end{aligned}
$$

The last equation suggests, roughly: the more spread out $f$ is, the more $\hat{f}$ will be concentrated at the origin and vice versa. This notion can be put in a precise form as follows.
(1.22) HEISENBERG INEQUALITY $(n=1:$ For $f \in S(\mathbb{R})$, we have $\|x f(x)\|_{2}\|\xi \hat{f}(\xi)\|_{2} \geq(1 / 4 \pi)\|f\|_{2}^{2}$.

Proof. Observe that

$$
\frac{d}{d x}(x f(x))=x \frac{d f}{d x}(x)+f(x)
$$

Thus

$$
\begin{aligned}
\|f\|_{2}^{2} & =\int f(x) \overline{f(x)} d x \\
& =\int \overline{f(x)}\left[\frac{d}{d x}(x f(x))-x \frac{d f}{d x}(x)\right] d x \\
& =-\int x f(x) \frac{d \bar{f}}{d x}(x) d x-\int x \frac{d f}{d x}(x) \bar{f}(x) d x \\
& =-2 \operatorname{Re} \int x f(x) \frac{d \bar{f}}{d x}(x) d x \\
& \leq 2\|x f(x)\|_{2}\left\|\frac{d \bar{f}}{d x}\right\|_{2} \text { (by Cauchy-Schwarz) } \\
\text { i.e., } \quad\|f\|_{2}^{2} & \leq 2\|x f(x)\|_{2}\left\|\frac{d f}{d x}\right\|_{2}
\end{aligned}
$$

But

$$
\left(\frac{d \bar{f}}{d x} \hat{)}(\xi)=2 \pi i \xi \hat{f}(\xi)\right.
$$

Therefore,

$$
\|f\|_{2}^{2} \leq 2.2 \pi\|x f(x)\|_{2}\|\xi \hat{f}(\xi)\|_{2}
$$

or

$$
\|x f(x)\|_{2}\|\xi f(\xi)\|_{2} \geq(1 / 4 \pi)\|f\|_{2}^{2}
$$

## A GENERALISATION TO $n$ VARIABLES

We replace $x$ by $x_{j}, \frac{d}{d x}$ by $\frac{\partial}{\partial x_{j}}$. Also for any $a_{j}, b_{j} \in \mathbb{R}$, we can replace $x_{j}$ and $\frac{\partial}{\partial x_{j}}$ by $x_{j}-a_{j}$ and $\frac{\partial}{\partial x_{j}}-b_{j}$ respectively. The same proof then yields:

$$
\begin{equation*}
\left\|\left(x_{j}-a_{j}\right) f(x)\right\|_{2}\left\|\left(\xi_{j}-b_{j}\right) \hat{f}(\xi)\right\|_{2} \geq(1 / 4 \pi)\|f\|_{2}^{2} . \tag{1.23}
\end{equation*}
$$

Let us now take $\|f\|_{2}=1$ and $A=\left(a_{1}, a_{2}, \ldots, a_{n}\right) \in \mathbb{R}^{n}$. Let $f$ be small outside a small nighbourhood of $A$.

In this case, $\left\|\left(x_{j}-a_{j}\right) f(x)\right\|_{2}$ will be small. Consequently, the other factor on the left in (1.23) has to be large. That is, if the mass of $f$ is concentrated near one point, the mass of $\hat{f}$ cannot be concentrated near any point.

Remark 1.24. If we take

$$
a_{j}=\int x_{i}|f(x)|^{2} d x, b_{j}=\int \xi|\hat{f}(\xi)|^{2} d \xi
$$

then inequality 1.23 is the mathematical formulation of the positionmomentum uncertainty relation in Quantum Mechanics.

## 3 Some Results From the Theory of Distributions

In this section, let us recall briefly some results from the theory of distributions. (For a more detailed treatment, see, for example [7] or [8]).

In the sequel, $\mathcal{D}^{\prime}(\Omega)$ will denote the space of distributions on the open set $\Omega \subset \mathbb{R}^{n}$ which is the dual space of $C_{o}^{\infty}(\Omega)$. When $\Omega=\mathbb{R}^{n}$, we will simply write $\mathcal{D}^{\prime}$ instead of $\mathcal{D}^{\prime}\left(\mathbb{R}^{n}\right)$. In the same way, $S^{\prime}=S^{\prime}\left(\mathbb{R}^{n}\right)$ will denote the space of tempered distributions with and $E^{\prime}=E^{\prime}\left(\mathbb{R}^{n}\right)$ will stand for the space of distributions with compact support.

The value of a distribution $u \epsilon \mathcal{D}^{\prime}$ at a function $\phi \epsilon C_{o}^{\infty}$ will be denoted by $\langle u, \phi\rangle$. If $u$ is a locally integrable function, then $u$ defines a distri-
write $\int u(x) \phi(x) d x$ for $\langle u, \phi\rangle$, when $u$ is an arbitrary distribution.
The convergence in $\mathcal{D}^{\prime}$ is the weak convergence defined by the following:
$u_{n}, u \epsilon \mathcal{D}^{\prime}, u_{n} \rightarrow$ in means $\left\langle u_{n}, \phi\right\rangle \rightarrow\langle u, \phi\rangle$ for every $\phi$ in $C_{o}^{\infty}$.
Let us now recall briefly certain operations on distributions.
(1.25) We can multiply a distribution $u \epsilon \mathcal{D}^{\prime}$ by a $C^{\infty}$ function $\phi$ to get another distribution $\phi u$ which is defined by

$$
<\phi u, \psi>=<u, \phi \psi>
$$

A $C^{\infty}$ function $\psi$ is said to be tempered if, for every multiindex $\alpha, \partial^{\alpha} \phi$ grows at most polynomially at $\infty$. We can multiply an $u, \epsilon S^{\prime}$ by a tempered function to get another tempered distribution. The definition is same as in the previous case.
(1.26) If $u \in \mathcal{D}^{\prime}$ and $f \in C_{o}^{\infty}$, we define the convolution $u * f$ by $(u * f)(x)=<u, f_{x}>$ where $f_{x}(x)=f(x-y)$.

The function $u * f$ is $C^{\infty}$ and when $u \in E, u * f$ is in $C_{o}^{\infty}$. The convolution $*: \mathcal{D}^{\prime} \times C_{o}^{\infty} \rightarrow C^{\infty}$ can be extended to a map from $\mathcal{D}^{\prime} \times E^{\prime}$ to $\mathcal{D}^{\prime}$. Namely, if $u \in \mathcal{D}^{\prime}, v \in E^{\prime}$ and $\phi \epsilon C_{o}^{\infty},\langle u * v, \phi\rangle=\langle u, \tilde{v} * \phi\rangle$ where $\tilde{v}$ is defined by $\langle\tilde{v}, \psi\rangle=\int v(x) \psi(-x) d x$. The associative law

$$
u *(v * w)-(u * w) * w \text { holds for } u, v, w \in \mathcal{D}^{\prime}
$$

provided that at most one of them does not have compact support. For $u \epsilon S^{\prime}$ and $f \epsilon S, u * f$ can also defined in the same way and $u * f$ is a tempered $C^{\infty}$ function.
(1.27) Since the Fourier transform is an isomorphism of $S$ onto $S$. and $\int f \hat{g}=\int \hat{f} g$, the Fourier transform extends by duality to an isomorphism of $S^{\prime}$ onto $S^{\prime}$.

For $u \epsilon E^{\prime} \subset S^{\prime}$, we have $\hat{u}(\xi)=<u, e^{-2 \pi i(\cdot) \cdot \xi}>$ which is an entire analytic function.

## Chapter 2

## Partial Differential Operators with Constant Coefficients

## 1 Local Solvability and Fundamental Solution

For the sake of convenience in taking the Fourier transforms, from now onwards, we will use the differential monomials

$$
D^{\alpha}=(2 \pi i)^{|\alpha|} \partial^{\alpha} . \text { Thus }\left(D^{\alpha} \hat{)}(\xi)=\xi^{\alpha} \hat{f}(\xi)\right.
$$

By a partial differential operator with constant coefficients, we mean a differential operator $L$ of the form

$$
L=\sum_{|\alpha| \leq k} a_{\alpha} D^{\alpha}, a_{\alpha} \epsilon \mathbb{C}
$$

Further, we assume that $\sum_{|\alpha| \leq k}\left|a_{\alpha}\right| \neq 0$. In this case, we say that the operator $L$ is of order $k$. If we write $P(\xi)=\sum_{|\alpha| \leq k} a_{\alpha} \xi^{\alpha}$, then we have $L=P(D)$.

For $u \in S$, taking the Fourier transform, we see that

$$
(P(D) u)(\xi)=\sum_{|\alpha| \leq k} a_{\alpha} \xi^{\alpha} \hat{u}(\xi)=P(\xi) \hat{u}(\xi)
$$

Let us now consider the following problem:
Given $f$ in $C^{\infty}$, we want to find a distribution $u$ such that $P(D) u=f$.
We say that the differential operator $L$ is locally solvable at $x_{0} \in \mathbb{R}^{n}$, if there is a solution of the above problem in some neighbourhood of the point $x_{\circ}$ for any $f$ in $C^{\infty}$.

17 Remark 2.1. We can assume that $f$ has compact support. To see this, we can take any $\phi$ in $C_{o}^{\infty}$ such that $\phi=1$ in some nighbourhood of the point $x_{0}$. If solve the problem $P(D) u=f \phi$ near $x_{0}$, then $u$ is a solution of our original problem since we have $f \phi=f$ in a neighbourhood of $x_{o}$.

In the following theorem we give an affirmative answer to the question of local solvability of $L=P(D)$. The simple proof exhibited here is due to $L$. Nirenberg.

Theorem 2.2. Let $L=\sum_{|\alpha| \leq k} a_{\alpha} D^{\alpha}$ be a differential operator with constant coefficients. If $u \in C_{o}^{\infty}$, there is a $C^{\infty}$ function $u$ satisfying $L u=f$ on $\mathbb{R}^{n}$.

Proof. Taking the Fourier transform of the equation $L u=P(D) u=f$, we see that $P(\xi) \hat{u}(\xi)=\hat{f}(\xi)$. It is natural to try to define $u$ by the formula

$$
u(x)=\int e^{2 \pi i x \cdot \xi} \hat{f}(\xi) / P(\xi) \xi
$$

In general, $P$ will have many zeros: so there will be a problem in applying the inverse Fourier transform to $\hat{f} / P$. But things are not so bad. Since $f \epsilon C_{o}^{\infty}, \hat{f}$ is an entire function of $\xi \in \mathbb{C}^{n}$ and $P$ is obviously entire. Hence we can deform the contour of integration to avoid the zeros of $P(\xi)$.

To make this precise, let us choose a unit vector $\eta$ so that $\sum_{|\alpha| \leq k} a_{\alpha} \eta^{\alpha} \neq$ 0 . By a rotation of coordinates, we can assume that $\eta=(0,0, \ldots, 0,1)$. Multiplying by a constant, we can also assume that $a_{\alpha 0}=1$ where $\alpha_{o}=(0,0, \ldots, 0, k)$. Then we have $P(\xi)=\xi_{n}^{k}+$ (lower order terms in $\left.\xi_{n}\right)$. Denote $\xi=\left(\xi^{\prime}, \xi_{n}\right)$ with $\xi^{\prime}=\left(\xi_{1}, \ldots, \xi_{n-1}\right)$ in $\mathbb{R}^{n-1}$.

Consider $P(\xi)=P\left(\xi^{\prime}, \xi_{n}\right)$ as a polynomial in the last variable $\xi_{n}$ in $\mathbb{C}$ for $\xi^{\prime}$ in $\mathbb{R}^{n-1}$. Let $\lambda_{1}\left(\xi^{\prime}\right), \ldots, \lambda_{k}\left(\xi^{\prime}\right)$ be the roots of $P\left(\xi^{\prime}, \xi_{n}\right)=0$ arranged so that if $i \leq j$.
$\mathfrak{J} \lambda_{i}\left(\xi^{\prime}\right) \leq \mathfrak{J} \lambda_{j}\left(\xi^{\prime}\right)$ and $\operatorname{Re} \lambda_{i}\left(\xi^{\prime}\right) \leq \operatorname{Re} \lambda_{i}\left(\xi^{\prime}\right)$ when $\mathfrak{J} \lambda_{i}\left(\xi^{\prime}\right)=\mathfrak{I} \lambda_{j}\left(\xi^{\prime}\right)$.
Since the roots of a polynomial depends continuously on the coefficients we see that $\lambda_{j}\left(\xi^{\prime}\right)$ are continuous in $\xi^{\prime}$. To proceed further, we need the help of two Lemmas.

Lemma 2.3. There is a measurable function $\phi: \mathbb{R}^{n-1} \rightarrow[-k-1, k+1]$ such that for all $\xi^{\prime}$ in $\mathbb{R}^{n-1} \min _{1 \leq j \leq k}\left\{\left|\phi\left(\xi^{\prime}\right)-\operatorname{Im} \lambda_{j}\left(\xi^{\prime}\right)\right|\right\} \geq 1$.

Proof. Left as an exercise to the reader: (cf. G.B. Folland [1]).
(The idea is that at least one of the $k+1$ intervals $[-k-1, k+1],[-k-$ $1, k+1] \ldots,[-k-1, k+1]$ must contain none of the $k$ points $\operatorname{Im} \lambda_{j}\left(\xi^{\prime}\right), j=$ $1,2, \ldots, k)$.

Lemma 2.4. Let $P(\xi)=\xi_{n}^{k}+($ lower order terms) and let $N(P)$ be the set $\left\{\zeta \epsilon \mathbb{C}^{n}: P(\zeta)=0\right\}$. Let $d(\xi, N(P))$. Then, we have

$$
|P(\xi)| \geq(d(\xi) / 2)^{k} .
$$

Proof. Take $\xi$ in $\mathbb{R}^{n}$ such that $P(\xi) \neq 0$. Let $\eta=(0,0, \ldots 0,1)$ and define $g(z)=P(\xi+z \eta)$ for $z$ in $\mathbb{C}$. This $g$ is a polynomial in one complex variable $z$. Let $\lambda_{1}, \lambda_{2}, \ldots, \lambda_{k}$ be the zeros of the polynomial $g$. Then

$$
g(z)=c\left(z-\lambda_{1}\right) \cdots\left(z-\lambda_{k}\right)
$$

so that

$$
\left|\frac{g(z)}{g(0)}\right|=\prod_{j=1}^{k}\left|1-\frac{z}{\lambda_{j}}\right| .
$$

Since $\xi+\lambda_{j} \eta \epsilon N(P)\left|\lambda_{j}\right| \geq d(\xi)$ so that when $|z| \leq d(\xi),\left|\frac{g(z)}{g(0)}\right| \leq 2^{k}$. Also

$$
\left|g^{(k)}(0)\right|=\left|\frac{k!}{2 \pi i} \int_{|\zeta|=d(\xi)} q(\zeta) \zeta^{-k-1} d \zeta\right| \leq \frac{k!}{2} \frac{|g(0)|}{d(\xi)^{k+1}} 2^{k} 2 \pi d(\xi)
$$

i.e. $\left|g^{(k)}(0)\right| \leq k!|g(0)| 2^{k} d(\xi)^{-k}$. But $g(0)=P(\xi)$ and

$$
\left|g^{(k)}(0)\right|=\frac{\partial^{k}}{\partial \xi_{n}^{k}} P(\xi)=k!
$$

Therefore,

$$
k!\leq k!|g(0)| 2^{k} d(\xi)^{-k}
$$

or

$$
|P(\xi)| \geq(d(\xi) / 2)^{k}
$$

Hence the lemma is proved.
Returning to the proof of the theorem, consider the function

$$
u(x)=\int_{\mathbb{R}^{n-1}} \int_{I M \xi_{n}=\phi\left(\xi^{\prime}\right)} e^{2 \phi i x \cdot \xi}(\hat{f}(\xi) / P(\xi)) d \xi_{n} d \xi
$$

By Lemmas 2.3 and 2.4 we have

$$
|P(\xi)| \geq(d(\xi) / 2)^{k} \geq 2^{-k} \text { along } \operatorname{Im} \xi_{n}=\phi\left(\xi^{\prime}\right)
$$

Since $f \epsilon C_{\circ}^{\infty}, \hat{f}(\xi)$ is rapidly decreasing as $|\operatorname{Re} \xi|$ tends to $\infty$ as along as $|\operatorname{Im} \xi|$ stays bounded, so the integral converges absolutely and uniformly together with all derivatives defining a $C^{\infty}$ function $u$.

Finally, by Cauchy's theorem,

$$
\begin{aligned}
P(D) u(x) & =\int_{\mathbb{R}^{n-1}} \int_{I M \xi_{n}=\phi\left(\xi^{\prime}\right)} e^{2 \phi i x \cdot \xi}(\hat{f}(\xi) / P(\xi)) d \xi_{n} d \xi^{\prime} \\
& =\int_{\mathbb{R}^{n}} e^{2 \phi i x \cdot \xi}(\hat{f}(\xi) d \xi=f(x) .
\end{aligned}
$$

This completes the proof of Theorem 2.2
Let us now consider the local solvability of $L=P(D)$ in the case when $f$ is a distribution.

As before, we remark that it suffices to take $f \epsilon E^{\prime}$. Further, it is enough to consider the case where $f=\delta$. Indeed, if $K$ satisfies $P(D) K=$ $\delta$, then we have for any $f \epsilon E^{\prime}, P(D)(K * f)=P(D) K * f=\delta * f=f$.

Definition 2.5. A distribution $K$ satisfying $P(D) K=\delta$ is called a fundamental solution or elementary solution of the differential operator $L=P(D)$.

A remarkable theorem due to MALGRANGE AND EHRENPREIS states that every differential operator $P(D)$ with constant coefficients has a fundamental solution. In fact, we can prove this result by a simple extension of the preceding argument.

Theorem 2.6. Every partial differential operator $P(D)$ with constant co-
efficients has a fundamental solution.
Proof. Proceeding as in the previous theorem, we try to define

$$
K(x)=\int_{\mathbb{R}^{n-1}} \int_{\operatorname{Im} \xi_{n}=\phi\left(\xi^{\prime}\right)} e^{2 \Pi i x \cdot \xi}(p(\xi))^{-1} d \xi_{n} d \xi^{\prime}
$$

Here, however, the integral may diverge at infinity. So, we consider the polynomial

$$
P_{N}(\xi)=P(\xi)\left(1+4 \Pi^{2} \sum_{j=1}^{n} \xi_{j}^{2}\right) N
$$

where $N$ is a large positive integer. Let

$$
K_{N}(x)=\int_{\mathbb{R}^{n-1}} \int_{\operatorname{Im\xi } \xi_{n}=\phi\left(\xi^{\prime}\right)} e^{2 \Pi i x \cdot \xi}\left(P_{N}(\xi)\right)^{-1} d \xi_{n} d \xi^{\prime}
$$

where $\phi$ is chosen appropriately for the polynomial $P_{N}$. Then on the region of integration, we have

$$
\left|P_{N}(\xi)\right| \geq c\left(1+|\xi|^{2}\right)^{N}
$$

so the integral will converge when $N>n / 2$. We claim that $P_{N}(D) K_{N}=$ $\delta$.

To see this, take $\phi \epsilon C_{o}^{\infty}$ and observe that the transpose of $P_{N}(D)$ is $P_{N}(-D)$. Thus

$$
<P_{N}(D) K_{N}, \phi>=<K_{N}, P_{N}(-D) \phi>
$$

$$
\begin{aligned}
& =\int_{\mathbb{R}^{n}} \int_{\mathbb{R}^{n-1}} \int_{I m \xi_{n}=\phi\left(\xi^{\prime}\right)} e^{2 \Pi i x \cdot \xi} \frac{\left(P_{N}(-D) \phi\right)(x)}{P_{N}(\xi)} d \xi_{n} d \xi^{\prime} d x \\
& =\int_{\mathbb{R}^{n-1}} \int_{I_{m}=\phi\left(\xi^{\prime}\right)}\left(P_{N}(\xi)\right)^{-1} d \xi_{n} d \xi^{\prime} \int_{\mathbb{R}^{n}} e^{2 \Pi i x \cdot \xi} P_{N}(-D) \phi(x) d x \\
& =\int_{\mathbb{R}^{n-1}} \int_{I m \xi_{n}=\phi\left(\xi^{\prime}\right)} \hat{\phi}(-\xi) d \xi_{n} d \xi^{\prime} \\
& =\int_{\mathbb{R}^{n}} \hat{\phi}(-\xi) d \xi=\phi(0)=<\delta, \phi>.
\end{aligned}
$$

Thus,

$$
\delta=P_{N}(D) K_{N}=P(D)(1-\Delta)^{N} K_{N}
$$

so if we put $K=(1-\cdot \Delta)^{N} K_{N}$ we have $P(D) K=\delta$. Thus $K$ is a fundamental solution of $P(D)$.

## 2 Regularity Properties of Differential Operators

Definition 2.7. The singular support of a distribution $f \epsilon D^{\prime}$ is defined to be the complement of the largest open set on which $f$ is a $C^{\infty}$ function. The singular support of $f$ will be denoted by sing $\operatorname{supp} f$.

Definition 2.8. Let $L=\sum_{|\alpha| \leq K} a_{\alpha}(x) D^{\alpha}$ where $a_{\alpha} \epsilon C^{\infty}$ be a differential operator. L is said to be hypoelliptic, if and only if, for any $u \in \mathcal{D}^{\prime}$, sing $\operatorname{supp} u \subset \operatorname{sing} \operatorname{supp} L u$. In other words, $L$ is hypoelliptic if and only if for any open set $\Omega \subset \mathbb{R}^{n}$ and any $u \in \mathcal{D}^{\prime}(\Omega)$, "LuєC $(\Omega)$ implies $u \epsilon C^{\infty}(\Omega)$ ".

Remark 2.9. The operator $L$ is said to be elliptic at $x \in \mathbb{R}^{n}$ if

$$
\sum_{|\alpha|=k} a_{\alpha}(x) \xi^{\alpha} \neq 0 \text { for every } \xi \text { in } \epsilon \mathbb{R}^{n} \backslash\{0\}
$$

$L$ is said to be elliptic if $L$ is elliptic at every $x \in \mathbb{R}^{n}$. Elliptic operators are hypoelliptic, as we shall prove later. This accounts for the name hypoelliptic.

We know that an ordinary differential operator with $C^{\infty}$ coefficients is hypoelliptic as long as the top order coefficient is non-zero. But this is not the case with partial differential operators as seen from the following

Example 2.10. Take the operator $L=\frac{\partial^{2}}{\partial x \partial y}$ in $\mathbb{R}^{2}$. The general solution of the equation $L u=0$ is given by $u(x, y)=f(x)+g(y)$ where $f$ and $g$ are arbitrary $C^{1}$ functions. This shows that $L$ is not hypoelliptic.

We observe that if $L$ is hypoelliptic, then every fundamental solution of $L$ is a $C^{\infty}$ function in $\mathbb{R}^{n} \backslash\{0\}$.

In the case of partial differential operators with constant coefficients, this is also sufficient for hypoellipticity. Indeed, we have the following

Theorem 2.11. Let $L$ be a partial differential operator with constant coefficients. Then the following are equivalent:
a) L is hypoelliptic.
b) Every fundamental solution of $L$ is $C^{\infty}$ in $\mathbb{R}^{2} \backslash\{0\}$.
c) At least one fundamental solution of $L$ is $C^{\infty}$ in $\mathbb{R}^{2} \backslash\{0\}$.

Proof. That (a) simple (b) follows from the above observation and that (b) implies (c) is completely trivial. The only nontrivial part we need to prove is that (c) implies (a). To prove this implication, we need

Lemma 2.12. Suppose $f \in \mathcal{D}^{\prime}$ is such that $f$ is $C^{\infty}$ in $\mathbb{R}^{n} \backslash\{0\}$ and $g \epsilon E^{\prime}$. Then $\operatorname{sing} \operatorname{supp}(f * g) \subset \operatorname{supp} g$.

Proof. Suppose $x \notin \operatorname{supp} g$. We will show that $f * g$ is $C^{\infty}$ in a neighbourhood of $x$.

Since $x \notin \operatorname{supp} g$, there exists an $\epsilon>0$ such that $B(x, \epsilon) \cap \operatorname{supp} g=\phi$. Choose $\phi \epsilon C_{0}^{\infty}(B(0, \epsilon / 2))$ such that $\phi=1$ on $B(0, \epsilon / 4)$. Now, $f * g=$ $(\phi f) * g+(1-\phi) f * g$. Since $(1-\phi) f$ is a $C^{\infty}$ function, $(1-\phi) f * g$ is $C^{\infty}$. Also

$$
\operatorname{Supp}(\phi f * g) \subset \operatorname{Supp} \phi f+\operatorname{Supp} g
$$

$$
\subset\{y: d(y, \operatorname{supp} g) \leq \epsilon / 2\}
$$

which does not intersect $B(x, \epsilon / 2)$. Therefore, in $B(x, \epsilon / 2), f * g=(1-$ $\phi) f * G$ which is $C^{\infty}$.

Proof of theorem 2.11. Let $K$ be a fundamental solution of $L$ such that $K$ is $C^{\infty}$ in $\mathbb{R}^{n} \backslash\{0\}$. Suppose $u \epsilon \mathcal{D}^{\prime}$ and $L u \epsilon C^{\infty}(\Omega)$ where $\Omega \subset \mathbb{R}^{n}$ is open.

For $x \epsilon \Omega$, pick $\epsilon>0$ small enough so that $\bar{B}(x, \epsilon) \subset \Omega$. Choose $\phi \epsilon C_{o}^{\infty}(B(x, \epsilon))$ so that $\phi=1$ on $B(x, \epsilon / 2)$. Then $L(\phi u)=\phi L u+v$ where $v=0$ on $B(x, \epsilon / 2)$ and also outside $B(x, \epsilon)$. We write

$$
K * L(\phi u)=K * \phi L u+K * v
$$

$\phi L u$ is a $C_{o}^{\infty}$ function so that $K * L u$ is a $C^{\infty}$ function. Also $K * v$ is a $C^{\infty}$ function on the ball $B(x, \epsilon / 2)$ by Lemma 2.12

Therefore $K * L(\phi u)$ is a $C^{\infty}$ function on $B(x, \epsilon / 2)$. But $K * L(\phi u)=$ $L K * \phi u=\delta * \phi u=\phi u$. Thus, $\phi u$ is a $C^{\infty}$ function on $B(x, \epsilon / 2)$. Since $\phi=1$ on $B(x, \epsilon / 2), u$ is a $C^{\infty}$ function on $B(x, \epsilon / 2)$. Since $x$ is arbitrary, this completes the proof.

## 3 Basic Operators in Mathematical Physics

In this section, we introduce the three basic operators in Mathematical Physics. In the following sections, we shall compute fundamental solutions for these operators and show how they can be applied to solve boundary value problems and to yield other information.
(i) LAPLACE OPERATOR: $\Delta=\sum_{j=1}^{n} \frac{\partial^{2}}{\partial x_{j}^{2}}$ in $\mathbb{R}^{n}$. If $u(x)$ represents electromagnetic potential (or gravitational potential ) and $\rho$ denotes the charge (resp. mass) density, then they are connected by the equation $\Delta u=-4 \pi \rho$. If the region contains no charge, i.e., if $\rho=0, \Delta u=0$. This equation is called the homogeneous Laplace equation, and its solutions are called harmonic functions.
(ii) HEAT OPERATOR: $L=\frac{\partial}{\partial t}-\Delta$ in $\mathbb{R}^{n+1}$. If $u(x, t)$ represents the temperature of a homogeneous body at the position $x$ and time $t$, then $u$ satisfies the heat equation $\frac{\partial u}{\partial t}-\Delta u=0$ in $\mathbb{R}^{n+1}$.
(iii) WAVE OPERATOR: $\square=\frac{\partial^{2}}{\partial t^{2}}-\Delta$ in $\mathbb{R}^{n+1}$. If $u(x, t)$ represents the amplitude of an electromagnetic wave in vacuum at position $x$ and time $t$, then $u$ satisfies the wave equation $\square u=0$ in $\mathbb{R}^{n+1}$.

The equation $\square u=0$ can also be used to describe other types of wave phenomena although in most cases, it is only an approximation valid for small amplitudes. More generally, the equation $\square u=f$ describes waves subject to a driving force $f$.

The Laplace operator $\Delta$ is an ingredient in all the above examples. The reason for this is that the basic laws of Physics are invariant under translation and rotation of coordinates which severely restricts the differential operators which can occur in them. Indeed, we have

Theorem 2.13. Let $L$ be a differential operator which is invariant under rotations and translations. Then there exists a polynomial $Q$ in one variable with constant coefficients such that $L=Q(\Delta)$.

Proof. Let $L=P(D)=\sum_{|\alpha| \leq K} a_{\alpha} D^{\alpha}$. Since $P(D)$ is invariant under translations, $a_{\alpha}$ are constants. Let $T$ be any rotation. We have

$$
\begin{aligned}
P(T \xi) \hat{\phi}(T \xi) & =(P(D) \phi)^{\wedge}(T \xi) \\
& =(P(D) \phi \circ T)^{\wedge}(\xi) \\
& =(P(D) \phi \circ T))^{\wedge}(\xi) \\
& =P(\xi) \hat{\phi}(T \xi), \text { for every } \phi .
\end{aligned}
$$

Thus $P(T \xi)=P(\xi)$ so that $P$ is rotation - invariant.
Write $P(\xi)=\sum_{j=0}^{k} P_{j}(\xi)$ where $P_{j}(\xi)$ is the part of $P$ which is homogeneous of degree $j$. We claim that each $P_{j}$ is rotation - invariant. Indeed, if $t$ is any real number,
$P(t \xi)=\sum_{j=0}^{k} t^{j} P_{j}(T \xi)$. For a rotation $T$, since $P(\xi)=P(T \xi)$, we have
so that

$$
\begin{aligned}
& P(t \xi)=\sum_{j=0}^{k} t^{j} P_{j}(T \xi) \\
& \sum_{j=0}^{k} t^{j}\left(P_{j}(T \xi)-P_{j}(\xi)\right)=0 .
\end{aligned}
$$

This being true for all $t$, it follows that $P_{j} \circ T=P_{j}$. But the only rotation-invariant functions which are homogeneous of degree $j$ are of the form $P_{j}(\xi)=c_{j}|\xi|^{j}$ where $c_{j}$ are constants.

Indeed, since $P_{j}$ is rotation-invariant, $P_{j}(\xi)$ depends only on $|\xi|$ and thus, for $|\xi| \neq 0$,

$$
P_{j}(\xi)=P_{j}\left(|\xi| \frac{\xi}{|\xi|}\right)=|\xi|^{j} P_{j}\left(\frac{\xi}{|\xi|}\right)=c_{j}|\xi|^{j}
$$

Since $|\xi|^{j}$ is not a polynomial when $j$ is odd, $c_{j}=0$ in that case. Therefore,

$$
P(\xi)=\sum c_{2} j|\xi|^{2 j}
$$

Taking $Q(x)=\sum\left(c_{2 j} /\left(-4 \pi^{2}\right)^{j}\right) x^{j}$, we get $P(D)=Q(\Delta)$.
Remark 2.14. This theorem applies to scalar differential operators. If one considers operators on vector or tensor functions, there are first order operators which are translation - and rotation - invariant, namely, the familiar operators grad, curl, div of 3-dimensional vector analysis and their $n$-dimensional generalisations.

Definition 2.15. A function $F(x)$ is said to be radial, if there is a function $f$ of one variable such that $F(x)=f(|x|)=f(r), r=|x|$.

When $F$ is radial and $F \epsilon L^{1}\left(\mathbb{R}^{n}\right)$, we have

$$
\int_{\mathbb{R}^{n}} F(x) d x=\int_{0}^{\infty} \int_{|x|=1} f(r) r^{n-1} d \sigma(x) d r
$$

where $d \sigma(x)$ is the surface measure on $S^{n-1}$, the unit sphere in $\mathbb{R}^{n}$.
Thus

$$
\int_{\mathbb{R}^{n}} F(x) d x=\omega_{n} \int_{0}^{\infty} f(r) r^{n-1} d r
$$

where $\omega_{n}$ is the area of $S^{n-1}$.
Let us now calculate $\omega_{n}$. We have

$$
\int e^{-\pi|x|^{2}} d x=1
$$

Now

$$
\begin{aligned}
\int e^{-\pi|x|^{2}} d x & =\omega_{n} \int_{0}^{\infty} e^{-\pi r^{2}} r^{n-1} d r \\
& =\frac{\omega_{n}}{2 \pi} \int_{0}^{\infty} e^{-s}(s / \pi)^{n / 2-1} d s, s=\pi r^{2} \\
& =\frac{\omega_{n}}{2 \pi^{n / 2}} \int_{0}^{\infty} e^{-s} s^{n / 2-1} d s \\
& =\omega_{n} \Gamma(n / 2) /\left(2 \pi^{n / 2}\right)
\end{aligned}
$$

Therefore

$$
\omega_{n}=2 \pi^{n / 2} / \Gamma(n / 2)
$$

## 4 Laplace Operator

First, let us find a fundamental solution of $\Delta$, i.e., we want to find a distribution $K$ such that $\Delta K=\delta$.

Since $\Delta$ commute with rotations and $\delta$ has the same property, we observe that if $u$ is a fundamental solution of $\Delta$ and $T$ is a rotation, then $u \circ T$ is also a fundamental solution. We therefore expect to find a fundamental solution $K$ which is radial.

Let us tarry a little to compute the Laplacian of a radial function $F$. $\operatorname{Let} F(x)=f(r), r=|x|$. Then,

$$
\begin{aligned}
\Delta F(x) & =\sum_{j=1}^{n} \frac{\partial}{\partial x_{j}}\left[f^{\prime}(r) x_{j} / r\right] \\
& =\sum_{j=1}^{n}\left(f^{\prime \prime}(r) \frac{x_{j}^{2}}{r^{2}}+\frac{f^{\prime}(r)}{r}-\frac{f^{\prime}(r)}{r^{3}} x_{j}^{2}\right) \\
& =f^{\prime \prime}(r)+((n-1) / r) f^{\prime}(r) .
\end{aligned}
$$

Now set $K(x)=f(r)$. If $K$ is to be a fundamental solution of $\Delta$, we must have

$$
f^{\prime \prime}(r)+((n-1) / r) f^{\prime}(r)=0 \text { on }(0, \infty) .
$$

From this equation

$$
f^{\prime \prime}(r) / f^{\prime \prime}(r)=-(n-1) / r
$$

Integrating, we get

$$
f^{\prime}(r)=c_{\circ} r^{1-n}
$$

with a constant $c_{0}$. One more integration yields

$$
\begin{aligned}
f(r) & =c_{1} r^{2-n}+c_{2} \text { when } n \neq 2 \\
& =c_{1} \log r+c_{2} \text { when } n=2 .
\end{aligned}
$$

Since constants are solutions of the homogeneous Laplace equation, we may assume that $c_{2}=0$. Thus if we set $F(x)=|x|^{2-n}(n \neq 2)$ or $F(x)=\log |x|(n=2)$, we expect to find that $\Delta F=c \delta$ for some $c \neq 0$, and then $K=c^{-1} F$ will be our fundamental solution.

In fact, we have
Theorem 2.16. If $F(x)=|x|^{2-n}$ on $\mathbb{R}^{n}(n \neq 2)$, then $\Delta F=(2-n) \omega_{n} \delta$, where $\omega_{n}$ is the area of the unit sphere in $\mathbb{R}^{n}$,

Proof. For $\epsilon>0$, we define $F_{\epsilon}(x)=\left(\epsilon^{2}+|x|^{2}\right)^{(2-n) / 2}$. Then $F_{\epsilon}$ is a $C^{\infty}$ function on $\mathbb{R}^{n}$, and an application of Lebesgue's Dominated Convergence Theorem reveals that $F_{\epsilon}$ converges to $F$ as $\epsilon$ tends to 0 , in the sense of distributions. Therefore $\Delta F_{\epsilon}$ converges to $\Delta F$, in the same sense. Let us now compute $\Delta F_{\epsilon}$.

$$
\begin{aligned}
\Delta F_{\epsilon}(x) & =\sum_{j=1}^{n} \frac{\partial}{\partial x_{j}}\left\{(2-n)\left(|x|^{2}+\epsilon^{2}\right)^{-n / 2} x_{j}\right\} \\
& =\sum_{j=1}^{n}\left\{(2-n)\left(|x|^{2}+\epsilon^{2}\right)^{-n / 2}+(2-n)(-n)\left(|x|^{2}+\epsilon^{2}\right)^{(-n / 2)-1} x_{j}^{2}\right\} \\
& =(2-n)\left(|x|^{2}+\epsilon^{2}\right)^{(-n / 2)-1} n \epsilon^{2} .
\end{aligned}
$$

Thus we see that $\Delta F_{\epsilon} \epsilon L^{1}\left(\mathbb{R}^{n}\right)$. A simple computation shows that $\Delta F_{\epsilon}(x)=\epsilon^{-n} \Delta F_{1}(x / \epsilon)$; so, by Theorem 1.6, $\Delta F_{\epsilon}$ tends to $\left(\int \Delta F_{1}\right) \delta$ as $\epsilon$ tends to 0 . Therefore $\Delta F=\left(\int \Delta F_{1}\right) \delta$, and we need only to compute $\int \Delta F_{1}$.

$$
\begin{aligned}
\int \Delta F_{1} & =n(2-n) \int\left(1+|x|^{2}\right)^{(-n / 2)-1} d x \\
& =n(2-n) \omega_{n} \int_{0}^{\infty}\left(1+r^{2}\right)^{(-n / 2)-1} r^{n-1} d r
\end{aligned}
$$

Putting $r^{2}+1=s$, we see that

$$
\begin{aligned}
\int \Delta F_{1} & =\frac{1}{2} n(2-n) \omega_{n} \int_{1}^{\infty} s^{(-n / 2)-1}(s-1)^{(n-1) / 2} d s \\
& =\frac{1}{2} n(2-n) \omega_{n} \int_{1}^{\infty} s^{-2}\left(1-\frac{1}{s}\right)^{(n / 2)-1} d s \\
& =\frac{1}{2} n(2-n) \omega_{n} \int_{0}^{\infty}(1-\sigma)^{(n / 2)-1} d \sigma, \sigma=\frac{1}{s} \\
& =\frac{1}{2} n(2-n) \omega_{n}(2 / n)=(2-n) \omega_{n}
\end{aligned}
$$

which completes the proof.
Exercise. Show that if $F(x)=\log |x|$ on $\mathbb{R}^{2}$, then $\Delta F=2 \pi \delta$.
Corollary 2.17. Let $K(x)=\left\{\begin{array}{l}\frac{|x|^{2-n}}{(2-n) \omega_{n}}(n \neq 2) \\ \frac{1}{2 \pi} \log |x|(n=2)\end{array}\right.$
Then $K$ is a fundamental solution of the Laplacian.
Corollary 2.18. $\Delta$ is hypoelliptic.
Proof. Follows from Theorem 2.11
It is also instructive to compute the fundamental solution of $\Delta$ by the Fourier transform method.

If $K$ is a fundamental solution of $\Delta$, we have $\Delta(k * f)=\Delta K * f=f$. Since $(\Delta g)^{\wedge}(\xi)=-4 \pi^{2}|\xi|^{2} \hat{g}(\xi)$, we get

$$
-4 \pi^{2}|\xi|^{2} \hat{K}(\xi) \hat{f}(\xi)=\hat{f}(\xi)
$$

so that, at least formally,

$$
\hat{K}(\xi)=-1 /\left(4 \pi^{2}|\xi|^{2}\right)
$$

We observe that when $n>2$, the function $-1 /\left(4 \pi^{2}|\xi|^{2}\right)$ is locally integrable and so defines a tempered distribution. We want to show that its inverse Fourier transform is our fundamental solution $K$. For that purpose, we will prove a more general theorem.
Theorem 2.19. For $0<\alpha<n$, let $F_{\alpha}$ be the locally integrable function $F_{\alpha}(\xi)=|\xi|^{-\alpha}$. Then

$$
F_{\alpha}^{\vee}(x)=\left(\Gamma\left(\frac{1}{2}(n-\alpha)\right) / \Gamma\left(\frac{1}{2} \alpha\right)\right) \pi^{\alpha-n / 2}|x|^{\alpha-n}
$$

Before proving this theorem, let us pause a minute to observe that it implies

$$
\begin{aligned}
F_{2}^{\vee}(x) & =\Gamma((n / 2)-1) \pi^{2-n / 2}|x|^{2-n}=(\Gamma(n / 2) /(n / 2-1)) \pi^{2-n / 2}|x|^{2-n} \\
& \left.=\left(-4 \pi^{2} /\left((2-n) \omega_{n}\right)\right)|x|^{2-n}, \text { so that }\left(-F_{2} / 4 \pi^{2}\right)\right)^{\vee}=K
\end{aligned}
$$

as desired.

Proof of theorem 2.19. The idea of the proof is to express $F_{\alpha}$ as a weighted average of Gaussian functions, whose Fourier transforms we can compute.

To begin with,

$$
\int_{0}^{\infty} e^{-r t} t^{\alpha-1} d t=\int_{0}^{\infty} e^{-s} s^{\alpha-1} r^{-\alpha} d s=r^{-\alpha} \Gamma(\alpha),(r>0, \alpha>0)
$$

In other words, for any $r>0$ and $\alpha>0$,

$$
r^{-\alpha}=\frac{1}{\Gamma(\alpha)} \int_{0}^{\infty} e^{-r t} t^{\alpha-1} d t
$$

Taking $r=\pi|\xi|^{2}$ and replacing $\alpha$ by $\alpha / 2$, then

$$
|\xi|^{-\alpha}=\frac{\pi^{\alpha / 2}}{\Gamma(\alpha / 2)} \int_{0}^{\infty} e^{-\pi|\xi|^{2 t} t} t^{(\alpha / 2)-1} d t
$$

which is the promised formula for $F_{\alpha}$ as a weighted average of Gaussians. Formally, we can write

$$
\begin{aligned}
\int_{\mathbb{R}^{n}} e^{2 \pi i x . \xi}|\xi|^{-\alpha} d \xi & =\frac{\pi^{\alpha / 2}}{\Gamma(\alpha / 2)} \int_{\mathbb{R}^{n}} \int_{0}^{\infty} e^{2 \pi i x \cdot \xi_{e}-\pi|\xi|^{2 t}} t^{\alpha / 2-1} d t d \xi \\
& =\frac{\pi^{\alpha / 2}}{\Gamma(\alpha / 2)} \int_{0}^{\infty}\left(e^{\pi|\xi|^{2 t}}\right)^{v}(x) t^{(\alpha / 2)-1} d t \\
& =\frac{\pi^{\alpha / 2}}{\Gamma(\alpha / 2)} \int_{0}^{\infty} e^{-\pi\left(|x|^{2}\right) / t} t^{-n / 2} t^{\alpha / 2-1} d t \\
& =\frac{\pi^{\alpha / 2}}{\Gamma(\alpha / 2)} \int_{0}^{\infty} e^{-\pi|x|^{2} S_{S}(n-\alpha) / 2-1} d s \\
& =\frac{\pi^{\alpha / 2}}{\Gamma(\alpha / 2)} \frac{\Gamma((n-\alpha) / 2)}{n^{(n-\alpha) / 2}}|x|^{\alpha-n} .
\end{aligned}
$$

In this computation, the change of order of the integration is unfortunately not justified, because the double integral is not absolutely convergent. This is not surprising, since $F_{\alpha}$ is not an $L^{1}$ function; instead, we must use the definition of the Fourier transform for distributions.

For every $\phi \epsilon S$,

$$
\begin{aligned}
<e^{-\pi|\xi|^{2 t}}, \hat{\phi}> & \left.=<e^{-\pi|\xi|^{2 t}}\right) \hat{,} \phi> \\
\text { i.e., } \quad \int_{\mathbb{R}^{n}} e^{-\pi|\xi|^{2 t}} \hat{\phi}(\xi) d \xi & =\int_{\mathbb{R}^{n}} t^{-n / 2} e^{-(\pi / t)|x|^{2}} \phi(x) d x .
\end{aligned}
$$

Now multiply both sides by $t^{(\alpha / 2)-1} \mathrm{dt}$ and integrate from 0 to $\infty$.

$$
\int_{0}^{\infty} \int_{\mathbb{R}^{n}} e^{-\pi|\xi|^{2 t}} t^{(\alpha / 2)-1} \hat{\phi}(\xi) d \xi d t=\int_{0}^{\infty} \int_{\mathbb{R}^{n}} e^{-\pi / t|x|^{2}} t^{(\alpha-n) / 2-1} \phi(x) d x d t
$$

The change of order in the integral is permitted now and integrating, we obtain,

$$
\begin{aligned}
\begin{aligned}
\frac{\Gamma(\alpha / 2)}{\pi^{\alpha / 2}} \int_{\mathbb{R}^{n}} \hat{\phi}(\xi)|\xi|^{-\alpha} d \xi & =\frac{\Gamma((n-\alpha) / 2)}{\pi^{(n-\alpha) / 2}} \int_{\mathbb{R}^{n}}|x|^{\alpha-n} \phi(x) d x \\
\text { i.e., } \quad<|\xi|^{-\alpha}, \hat{\phi}> & =<\frac{\Gamma((n-\alpha) / 2)}{(\Gamma(\alpha / 2)} \pi^{\alpha-n / 2}|x|^{\alpha-n}, \phi>.
\end{aligned}
\end{aligned}
$$

This shows that

$$
F_{\alpha}^{\vee}(x)=\frac{\Gamma((n-\alpha) / 2)}{\Gamma(\alpha / 2)} \pi^{\alpha-n / 2}|x|^{\alpha-n}
$$

as desired.
This analysis does not suffice to explain $-\left(4 \pi^{2}|\xi|^{2}\right)^{-1}$ as (in some sense) the Fourier transform of $(2 \pi)^{-1} \log |x|$ in the case $n=2$. One way to proceed is to define $G_{\alpha}(\xi)=(2 \pi|\xi|)^{-\alpha}$ and to study the behaviour of $G_{\alpha}$ and $G_{\alpha}^{\vee}$ as $\alpha$ tends to 2 . This analysis can be carried out just as easily in $n$ dimensions where the problem is to study $G_{\alpha}$ and $G_{\alpha}^{v}$ as $\alpha$ tends to $n$.

Let $R_{\alpha}=G_{\alpha}^{\vee}$ for $0<\alpha<n$. ( $R_{\alpha}$ called the Riesz potential of order $\alpha$ ). By the preceding theorem, we have

$$
R_{\alpha}(x)=\frac{(\Gamma((n-\alpha) / 2)}{2^{\alpha} \pi^{n / 2} \Gamma(\alpha / 2)}|x|^{\alpha-n}
$$

Moreover, if $f \epsilon S$,

$$
(-\Delta)^{\alpha / 2}\left(f * R_{\alpha}\right)=f
$$

in the sense that $(2 \pi|\xi|)^{\alpha}\left(f * R_{\alpha}\right)^{\wedge}(\xi)=\hat{f}(\xi)$.
As $\alpha$ tends to $n$, the $\Gamma$-function in $R_{\alpha}$ blows up. However, since $(2 \pi|\xi|)^{\alpha_{\delta}}=0$, we see that $(-\Delta)^{\alpha / 2} 1=0$; so we can replace $R_{\alpha}$ by $R_{\alpha}-c, c$ being a constant, still having

$$
(-\Delta)^{\alpha / 2}\left(f *\left(R_{\alpha}-c\right)\right)=f
$$

If we choose $c=c_{\alpha}$ appropriately, we can arrange that $R_{\alpha}-c_{\alpha}$ will have a limit as $\alpha$ tends to $n$. In fact, let us take

$$
c_{\alpha}=\frac{\Gamma((n-\alpha) / 2)}{2^{\alpha} \pi^{n / 2} \Gamma(\alpha / 2)} \text { and define } R_{\alpha}^{\prime}=R_{\alpha}-c_{\alpha} .
$$

Then

$$
\begin{aligned}
R_{\alpha}^{\prime}(x) & =\frac{\Gamma((n-\alpha) / 2)}{2^{\alpha} n^{n / 2} \Gamma(a / 2)}\left(|x|^{\alpha-n}-1\right) \\
& =\frac{2 \Gamma((n-\alpha) / 2)+1}{2^{\alpha} \pi^{n / 2} \Gamma(\alpha / 2)} \frac{(|x|)^{\alpha-n}-1}{n-\alpha} .
\end{aligned}
$$

Letting $\alpha$ tend to $n$, we get in the limiting case

$$
R_{n}^{\prime}(x)=\frac{-2^{1-n}}{\pi^{n / 2} \Gamma(n / 2)} \log |x|
$$

when $n=2,-R_{2}^{\prime}(x) \frac{1}{2 \pi} \log |x|$ and $\Delta\left(f *\left(-R_{2}^{\prime}\right)\right)=f$, so we recover our fundamental solution for the case $n=2$.

It remains to relate the function $G_{n}(\xi)=(2 \pi|\xi|)^{-n}$ to the Fourier 36
transform of the tempered distribution $F_{n}^{\prime}$. One way to make $G_{n}$ into a distribution is the following.

Define a functional $F$ on $S$ by

$$
<F, \phi>=\int_{|\xi| \leq 1} \frac{\phi(\xi)-\phi(0)}{(2 \pi|\xi|)^{n}} d \xi+\int_{|\xi|>1} \frac{\phi(\xi)}{(2 \pi|\xi|)^{n}} d \xi .
$$

Note that the first integral converges, since, in view of the mean value theorem, $|\phi(\xi)-\phi(0)| \leq c|\xi|$. It is easy to see that $F$ is indeed a tempered distribution. Further, we observe that when $\phi(0)=0$, $<F, \phi>=\int \phi(\xi) G_{n}(\xi) d \xi$, i.e., $F$ agrees with $G_{n}$ on $\mathbb{R}^{n} \backslash\{0\}$.

Just as $R_{n}^{\prime}$ was obtained from $R_{n}$ by subtracting off an infinite constant, $F$ is obtained from $G_{n}$ by subtracting off an infinite multiple of the $\delta$ function. This suggests that $F$ is essentially the Fourier transform of $R_{n}^{\prime}$. In fact, one has the following

Exercise. Show that $F-\left(R_{n}^{\prime}\right)^{\wedge}$ is a multiple of $\delta$.
We now examine a few of the basic properties of harmonic functions, i.e., functions satisfying $\Delta u=0$. We shall need the following results from advanced calculus.

Theorem 2.20 (DIVERGENCE THEOREM). Let $\Omega$ be a bounded open set in $\mathbb{R}^{n}$ with smooth boundary $\partial \Omega$. Let $v$ be the unit outward normal vector on $\partial \Omega$. Let $F: \mathbb{R}^{n} \rightarrow \mathbb{R}^{n}$ be a $C^{1}$ function on $\bar{\Omega}$, the closure of $\Omega$. Then we have

$$
\int_{\partial \Omega} F . v d \sigma=\int_{\Omega} d i v F d x=\int_{\Omega} \sum_{j=1}^{n} \frac{\partial F_{j}}{\partial x_{j}} d x .
$$

## CONSEQUENCES OF THEOREM 2.20

(2.21) When we take $F=\operatorname{grad} u, F . v=\operatorname{grad} u . v=\frac{\partial u}{\partial v}$, the normal derivative of $u$, and div $F=\operatorname{div} \operatorname{grad} u=\Delta u$. Therefore

$$
\int_{\partial \Omega} \frac{\partial u}{\partial v} d \sigma=\int_{\Omega} \Delta u d x
$$

(2.22) When we take $F=u \operatorname{grad} v-v \operatorname{grad} u$, div $F=u \Delta v-v \Delta u$. Therefore

$$
\int_{\partial \Omega}\left(u \frac{\partial v}{\partial v}-v \frac{\partial u}{\partial v}\right) d \sigma=\int_{\Omega}(u \Delta v-v \Delta u) d x
$$

This formula is known as Green's formula.
For harmonic functions, we have the following mean value theorem.
Theorem 2.23. Let $u$ be harmonic in $B\left(x_{0}, r\right)$. Then,for any $\rho<r$, we have

$$
u\left(x_{0}\right)=\frac{1}{\omega_{n} \rho^{n-1}} \int_{\left|x-x_{0}\right|=\rho} u(x) d \sigma(x) .
$$

That is, $u\left(x_{0}\right)$ is the mean value of $u$ on any sphere centred at $x_{0}$
Proof. Without loss of generality, assume that $x_{0}=0$. Since $\Delta u=$ 0 and $\Delta$ is hypoelliptic, $u$ is $a C^{\infty}$ function. Now formally,

$$
\begin{aligned}
u(0) & =<\delta, u> \\
& =\int_{|x|<\rho} \delta(x) u(x) d x \\
& =\int_{|x|<\rho} \Delta K(x) u(x) d x \\
& =\int_{|x|<\rho}(u \Delta K-K \Delta u) d x \\
& =\int_{|x|=\rho}\left(u \frac{\partial K}{\partial v}-K \frac{\partial u}{\partial v}\right) d \sigma, \text { by }(2.22) .
\end{aligned}
$$

Of course, we are cheating here by applying Green's formula to the non-smooth function $K$. Nonetheless the result is correct, and we leave it as an exercise for the reader to justify it rigorously (either approximate $K$ by $C^{\infty}$ functions as in the proof of Theorem 2.16 or apply Green's formula to the region $\epsilon<|x|<\rho$ and let $\epsilon$ tend to 0 .)

On the circle $|x|=\rho, K$ is a constant and so $\int_{|x|=\rho} K \frac{\partial u}{\partial v} d \sigma=$ Const $\int_{|x|=\rho} \frac{\partial u}{\partial v} d \sigma=$ Const $. \int_{|x| \leq \rho} \Delta u d \sigma=0$.

Further on $|x|=\rho \cdot \frac{\partial}{\partial v}=\frac{\partial}{\partial r}$, the radial derivative. Therefore

$$
\frac{\partial K}{\partial v}=1 /\left(\omega_{n} \rho^{n-1}\right) \text { on }|x|=\rho
$$

Thus we have

$$
u(0)=\frac{1}{\omega_{n} \rho^{n-1}} \int_{|x|=\rho} u(x) d \sigma \quad \text { as desired }
$$

As a corollary to the mean value theorem, we can derive the maximum principle for harmonic functions.

39 Theorem 2.24. Suppose $\Omega$ is a connected open set in $\mathbb{R}^{n}$. Let $u$ be a realvalued function harmonic in $\Omega$. If $A=\sup _{\Omega} u(x)$, then either $u(x)<A$ for all $x \in \Omega$ or $u(x) \equiv A$ on $\Omega$.

Proof. Suppose that $u\left(x_{0}\right)=A$ for some $x_{0} \epsilon \Omega$. From the mean value theorem

$$
u\left(x_{0}\right)=\frac{1}{\omega_{n} \rho^{n-1}} \int_{\left|x-x_{0}\right|=\rho} u(x) d \sigma(x)
$$

where $\rho$ is small enough so that $\left\{x:\left|x-x_{o}\right| \leq \rho\right\} \subset \Omega$. By our assumption, the integral does not exceed $A$. If $u\left(x_{1}\right)<A$ for some point of $B\left(x_{0}, \rho\right)$ then $u(x)<A$ in some neighbourhood of $x_{1}$ by continuity. Taking $r=\left|x_{1}-x_{0}\right|<\rho$, we have

$$
u\left(x_{0}\right)=\frac{1}{\omega_{n} r^{n-1}} \int_{\left|x-x_{o}\right|=r} u(x) d \sigma(x)<A
$$

a contradiction. Therefore, if we set $\Omega_{1}=\{x \in \Omega: u(x)=A\}$, then $\Omega_{1}$ is an open subset of $\Omega$ and $\Omega_{1} \neq \phi$. Further $\Omega_{2}=\{x \in \Omega: u(x)<A\}$ is also on open subset of $\Omega$ and $\Omega_{1} \cup \Omega_{2}=\Omega$. The connectedness of $\Omega$ and the non-emptiness of $\Omega_{1}$ force $\Omega_{2}$ to be empty. Thus $u(x) \equiv A$ on $\Omega$.

Corollary 2.25. If $\Omega$ is a bounded open set in $\mathbb{R}^{n}, u \in C(\bar{\Omega})$ and $\Delta u=$ 0 in $\Omega$, then

$$
\sup _{\Omega}|u(x)|=\sup _{\partial \Omega}|u(x)| .
$$

Proof. Since the function $u$ is continuous on the compact set $\bar{\Omega},|u(x)|$ attains its supremum on $\bar{\Omega}$ at some point $x_{0}$. By multiplying $u$ by a constant, we may assume that $u\left(x_{0}\right)=\left|u\left(x_{0}\right)\right|$.

If $x_{0} \epsilon \partial \Omega$, the corollary is proved. Otherwise $x_{0} \epsilon \Omega$ and by the previous theorem, on the connected component of containing $x_{0}, \operatorname{Re} u(x)=$ $u\left(x_{0}\right)$ and hence $\operatorname{Im} u(x)=0$. Since $u$ is continuous on $\bar{\Omega}, u(x)=u\left(x_{0}\right)$ for all $x \in \partial \Omega$. Thus $\sup _{\Omega}|u(x)|=\sup _{\partial \Omega}|u(x)|$.

Corollary 2.26. If $u$ and $v$ are in $C(\bar{\Omega}), \Omega$ is bounded $\Delta u=\Delta v=0$ on $\Omega$ and $u=v$ on $\partial \Omega$, then $u=v$ everywhere.

Proof. Apply the previous corollary to $u-v$.
The following boundary value problem for Laplace's equation, known as the Dirichlet problem, is of fundamental importance: given a function $f$ on $\partial \Omega$, find a function $u$ such that $\Delta u=0$ on $\Omega$ and $u=$ $f$ on $\partial \Omega$.

When $\Omega$ is bounded and $f \epsilon C(\partial \Omega)$ the uniqueness of the solution, if it exists at all, is assured by corollary 2.26. The problem of proving the existence of solution of the Dirichlet problem is highly nontrivial. We shall solve a special case, namely, when $\Omega$ is a half-space and then indicate how similar ideas can be applied for other regions.

First of all a word about notation: we will replace $\mathbb{R}^{n}$ by $\mathbb{R}^{n+1}$ with coordinates $\left(x_{1}, x_{2}, \ldots, x_{n}, t\right)$. Our Laplacian in $\mathbb{R}^{n+1}$ will be denoted by $\partial_{t}^{2}+\Delta$ where $\partial_{t}=\frac{\partial}{\partial t}$ and $\Delta=\sum_{j=1}^{n} \frac{\partial^{2}}{\partial x_{j}}$. Now our Dirichlet Problem is the following :
given a function $f$ on $\mathbb{R}^{n}$, find a function $u$ such that

$$
\begin{align*}
\left(\partial_{t}^{2}+\Delta\right) u(x, t) & =0, x \in \mathbb{R}^{n}, t>0  \tag{2.27}\\
u(x, o) & =f(x), x \in \mathbb{R}^{n}
\end{align*}
$$

Since the half space is unbounded, the uniqueness argument given above does not apply. Indeed, without some assumption on the behaviour of $u$ at $\infty$, uniqueness does not hold. For example, if $u(x, t)$ is a solution, so is $u(x, t)+t$. However, we have

Theorem 2.28. Let $u$ be a continuous function on $\mathbb{R}^{n} \times[0, \infty)$ satisfying
i) $\left(\Delta+\partial_{t}^{2}\right) u=0$ on $\mathbb{R}^{n} \times(0, \infty)$,
ii) $u(x, o)=0$ for $x \in \mathbb{R}^{n}$, and
iii) $u$ vanishes at $\infty$.

Then $U \equiv 0$ on $\mathbb{R}^{n} \times[0, \infty)$.
Proof. Applying the maximum principle for $u$ on $B(0, R) \times(0, T)$, we see that maximum of $|u|$ on $B(0, R) \times(0, T)$ tends to zero as $R, T$ tend $\infty$. Hence $u \equiv 0$.

Remark 2.29. Hypothesis (iii) in Theorem 2.28 can be replaced by (iii): $u$ is bounded on $\mathbb{R}^{n} \times[0, \infty)$. But this requires a deeper argument (See Folland [1]).

To solve the Dirichelt problem, we apply the Fourier transform in the variable $x$. We denote by $\tilde{u}(\xi, t)$ the Fourier transform

$$
\tilde{u}(\xi, t)=\int e^{-2 \pi i x \cdot \xi} u(x, t) d x
$$

If we take Fourier transform of (2.27), we obtain

$$
\begin{aligned}
\left(\partial_{t}^{2}-4 \pi^{2}|\xi|^{2}\right) \tilde{u}(\xi, t) & =o \text { on } \mathbb{R}^{n} \times(0, \infty) \\
\tilde{u}(\xi, o) & =\hat{f}(\xi) \text { on } \mathbb{R}^{n}
\end{aligned}
$$

The general solution of the ordinary differential equation (ODE) $\left(\partial_{t}^{2}-4 \pi^{2}|\xi|^{2}\right) \tilde{u}(\xi, t)=0$ is given by $\tilde{u}(\xi, t)=a(\xi) e^{-2 \pi|\xi| t}+b(\xi) e^{2 \pi|\xi| t}$ and the initial condition is $a(\xi)+b(\xi)=\hat{f}(\xi)$.

This formula for $\tilde{u}$ will define a tempered distribution in $\xi$, provided that as $|\xi|$ tends to $\infty,|a(\xi)|$ grows at most polynomially and $|b(\xi)|$
decreases faster than exponentially. We remove the non- uniqueness by requiring that $u$ should satisfy good estimates in terms of $f$ which are uniform in t -for example, that $u$ should be bounded if $f$ is bounded. This clearly forces $b(\xi)=0$, since $e^{2 \pi|\xi| t}$ blows up as $t$ tends to $\infty$. Thus we take $\tilde{u}(\xi, t)=\hat{f}(\xi) e^{-2 \pi|\xi|} t$. Taking the inverse Fourier transform in the variable $\xi, u(x, t)=\left(f * P_{t}\right)(x)$ where $P_{t}(x)=\left[e^{2 \pi t(|.|)}\right]^{\vee}(x)$ is the Poisson Kernel for $\mathbb{R}^{n} \times[0, \infty)$.

We now compute $P_{t}$ explicitly. When $n=1$, this is easy:

$$
\begin{aligned}
P_{t}(x) & =\int_{-\infty}^{0} e^{2 \pi i x \xi} e^{2 \pi \xi t} d \xi+\int_{0}^{\infty} e^{2 \pi i x \xi} e^{-2 \pi \xi t} d \xi \\
& =\frac{1}{2 \pi}\left[(t+i x)^{-1}+(t-i x)^{-1}\right]=\frac{t}{\pi}\left(x^{2}+t^{2}\right)^{-1}
\end{aligned}
$$

To compute $P_{t}(x)$ for the case $n>1$, as in the proof of Theorem 2.19, the idea is to express $e^{-2 \pi|\xi| t}$ as a weighted average of Gaussians.

Here's how!
Lemma 2.30. When $\beta>0, e^{-\beta}=\int_{0}^{\infty} \frac{e^{-s} e^{-\beta^{2} / 4 s}}{\sqrt{ }(\pi s)} d s$.
Proof. First, observe that

$$
e^{-\beta}=\frac{1}{\pi} \int_{-\infty}^{\infty} \frac{e^{i \beta \tau}}{\tau^{2}+1} d \tau
$$

This can be proved by using the residue theorem or by applying the Fourier inversion theorem to $P_{1}(x)=\frac{1}{\pi}\left(1+x^{2}\right)^{-1}$ on $\mathbb{R}^{1}$. We also have

$$
\frac{1}{\tau^{2}+1}=\int_{0}^{\infty} e^{-\left(\tau^{2}+1\right) s} d s
$$

Therefore,

$$
e^{-\beta^{\prime}}=\frac{1}{\pi} \int_{-\infty}^{\infty} \int_{0}^{\infty} e^{i 2 \pi \tau} e^{-\left(\tau^{2}+1\right) s} d s d \tau
$$

Putting $\tau=2 \pi \sigma$ and changing the order of integration,

$$
\begin{aligned}
& e^{=\beta}=\frac{1}{\pi} \int_{0}^{\infty} \int_{-\infty}^{\infty} e^{i 2 \pi \sigma \beta} e^{-s} e^{-4 \pi^{2} \sigma^{2} s} d \sigma d s \\
& \text { i.e., } \\
& e^{-\beta}=2 \int_{0}^{\infty} e^{-2}(4 \pi s)^{-1 / 2} e^{\beta^{2} / 4 s} d s \\
& =\int_{0}^{\infty} \frac{e^{-s} e^{\beta^{2} / 4 s}}{\sqrt{ }(\pi s)} d s
\end{aligned}
$$

which is the required expression.
Returning to the computation of $P_{t}(x)$, we have

$$
P_{t}(x)=\int e^{2 \pi i x \cdot \xi} e^{-2 \pi|\xi| t} d \xi
$$

Taking $\beta=2 \pi|\xi| t$ in the lemma, we get

$$
\begin{aligned}
P_{t}(x) & =\iint_{0}^{\infty} e^{2 \pi i x \cdot \xi} \frac{e^{-s}}{\sqrt{\pi s}} e^{\left(4 \pi^{2}|\xi|^{2} t^{2}\right) / 4 s} d s d \xi \\
& =\int_{0}^{\infty} \frac{e^{-s}}{\sqrt{\pi s}} \int e^{2 \pi i x \xi} e^{-\left(\pi^{2}|\xi|^{2} t^{2}\right) / S} d \xi \cdot d s \\
& =\int_{0}^{\infty} \frac{e^{-s}}{\sqrt{\pi s}}\left(\frac{\pi t^{2}}{s}\right)-n / 2 e^{-\left(s|x|^{2}\right) / t^{2}} d s \\
& =\pi^{-(n+1) / 2} t^{-n} \int_{0}^{\infty} e^{-s\left(1+\left(|x|^{2} / t^{2}\right)\right)} s^{(n-1) / 2} d s
\end{aligned}
$$

$$
=\frac{\pi^{-(n+1) / 2} t^{-n} \Gamma(n+1) / 2}{\left(1+\frac{|x|^{2}}{t^{2}}\right)(n+1 / 2)} .
$$

Thus we have

$$
\begin{equation*}
P_{t}(x)=\Gamma((n+1) / 2) \pi^{(n+1) / 2} t /\left(t^{2}+|x|^{2}\right)^{(n+1) / 2} \tag{2.31}
\end{equation*}
$$

In particular, we see that $P_{t} \epsilon L^{1} \cap L^{\infty}$, so that $f * P_{t}$ is well defined if, for example, $f \epsilon L^{p}, 1 \leq p \leq \infty$. That $\left(\Delta+\partial_{t}^{2}\right) P_{t}=0$ follows by taking the Fourier transform, and hence

$$
\left(\Delta+\partial_{t}^{2}\right)\left(f * P_{t}\right)=f *\left(\Delta+\partial_{t}^{2}\right) P_{t}=0 \text { for any } f
$$

Further note that

$$
P_{t}(x)=t^{-n} P_{1}(x / t) \text { and } \int P_{1}(x) d x=\hat{P}_{1}(0)=1
$$

Therefore, by Theorem 1.6, when $f \epsilon L^{p}, 1 \leq p<\infty, f * P_{t}$ tends to $f$ in the $L^{p}$ norm and when $f$ is bounded and continuous $f * P_{t}$ tends to $f$ uniformly on compact sets as $t$ tends to 0 .

If we take $f$ continuous and bounded, then for $u(x, t)=P_{t} * f$, we have $\lim _{t \rightarrow 0} u(x, t)=u(x, 0)=f(x)$. Thus the function $u(x, t)=\left(f * P_{t}\right)(x) \quad 45$ is a solution of the Dirichlet Problem for the half space.

Remark 2.32. The Poisson Kernel $P_{t}$ is closely related to the fundamental solution $K_{n+1}$ of the Laplacian in $\mathbb{R}^{n+1}$. Indeed

$$
K_{n+1}(x, t)=\frac{1}{(1-n) \omega_{n+1}}\left(|x|^{2}+t^{2}\right)^{-((n-1) / 2)}
$$

and hence

$$
P_{t}(x)=2 \partial_{t} K_{n+1}(x, t) .
$$

Exercise. Check the above equation using the Fourier transform. (Start with $\left(\partial_{t}^{2}+\Delta\right) K_{n+1}=\delta(x) \delta(t)$ and take the Fourier transform in the variable x to obtain

$$
\left(\partial_{t}^{2}-4 \pi^{2}|\xi|^{2}\right) \tilde{K}_{n+1}=\delta(t)
$$

Solve this equation to obtain

$$
\tilde{K}_{n+1}(\xi, t)=\frac{e^{2 \pi}|\xi \|||t|}{4 \pi|\xi|}
$$

Then

$$
\left.\partial_{t} \tilde{K}_{n+1}=\frac{1}{2} e^{-2 \pi|\xi| t}, t>0, \text { so that } \partial_{t} \tilde{K}_{n+1}=\frac{1}{2} \hat{P}_{t}\right)
$$

Our formula (2.31) for $P_{t}$ makes sense even when $t<0$ and we have $P_{-t}(x)=-P_{t}(x)$, so that $\lim _{t \rightarrow o \pm} f * P_{t}= \pm f$. We further observe that

$$
f *_{(x)} P_{t}=f(x) \delta(t) *_{(x, t)^{2}} \partial_{t} K_{n+1}(x, t)
$$

where $\left.*_{( } x\right)$ and $*_{(x, t)}$ mean convolution on $\mathbb{R}^{n}$ and $\mathbb{R}^{n+1}$ respectively.
Therefore,

$$
\begin{array}{ll} 
& f *_{(x)} P_{t}=f(x) \delta^{\prime}(t) *(x, t) 2 K_{n+1}(x, t) \\
\text { i.e., } \quad u(x, t)=2 f(x) \delta^{\prime}(t) *(x, t) K_{n+1}(x, t) .
\end{array}
$$

Form this, we

$$
\left(\delta_{t}^{2}+\Delta\right) u(x, t)=2 f(x) \delta^{\prime}(t) * \delta(x) \delta(t)=2 f(x) \delta^{\prime}(t)
$$

Exercise. Show directly that if
i) $\left(\Delta+\partial_{t}^{2}\right) u=0$ on $\mathbb{R}^{n+1} \backslash \mathbb{R}^{n} \times\{0\}$
ii) $u(x,-t)=-u(x, t)$
iii) $\lim _{t \rightarrow 0 \pm} u(x, t)= \pm f(x)$, then $u$ is a distribution solution of $\left(\partial_{t}^{2}+\Delta\right) u=$ $2 f(x) \delta^{\prime}(t)$.
[To avoid technicalities, assume $f$ is sufficiently smooth so that $\lim _{t \rightarrow 0 \pm} \frac{\partial u}{\partial t}$ exists; e.g. $f \epsilon C^{2}$ is sufficient ].

We now indicate, without giving any proof, how these ideas can be used to solve the Dirichlet Problem in a bounded open set $\Omega$ in $\mathbb{R}^{2}$. We assume that the boundary $\partial \Omega$ of $\Omega$ is of class $C^{2}$.

We know that the solution of the Dirichlet Problem in the case when $\Omega=\mathbb{R}^{n-1} \times(0, \infty)$ is given by $u=f * 2 \frac{\partial K}{\partial x_{n}}$ (the convolution is in $\mathbb{R}^{n-1}$ ). we note that $\frac{\partial K}{\partial x_{n}}$ is the inward normal derivative of $K$. For bounded $\Omega$ with $C^{2}$ boundary let us consider

$$
u(x)=2 \int_{\partial \Omega} f(y) \frac{\partial}{\partial v_{y}} K(x-y) d \sigma(y)
$$

Here $\sigma(y)$ is the surface measure on the boundary, $v$ is the unit outward normal on $\partial \Omega$, and $\frac{\partial \phi}{\partial v_{y}}(x, y)=\Sigma v_{j} \frac{\partial \phi}{\partial y_{j}}$. The minus sign in $K(x-y)$ compensates the switch from inward to outward normal.

Since $\Delta K(x-y)=\delta(x-y)$ we see that $\Delta u=0$ in $\Omega$. What about the behaviour of $u$ on $\partial \Omega$ ? It turns out that if $u$ is defined as above, then

$$
u \epsilon C(\bar{\Omega}) \text { and }\left.u\right|_{\partial \Omega}=f+T f
$$

where $T$ is a compact integral operator on $L^{2}(\partial \Omega)$ or $C(\partial \Omega)$. Hence if we can find $\phi$ in $C(\partial \Omega)$ such that $\phi+T \phi=f$, and we set

$$
u(x)=2 \int_{\partial \Omega} \phi(y) \frac{\partial}{\partial v_{y}} K(x-y) d \sigma(y)
$$

then $u$ satisfies $\Delta u=0$ in $\Omega$ and further $\left.u\right|_{\partial \Omega}=\phi+T \phi=f$. Hence the Dirichlet problem is reduced to solving the equation $\phi+T \phi=f$, and for this purpose, the classical Fredholm-Riesz theory is available. The upshot is that a solution to the Dirichlet problem always exists (provided, as we have assumed, that $\partial \Omega$ is of class $C^{2}$ ). See Folland [1] for a detailed treatment.

## 5 The Heat Operator

The Heat operator is given by $\partial_{t}-\Delta$. We want to find a distribution $K$ such that $\left(\partial_{t}-\Delta\right) K=\delta(x) \delta(t)$. Taking the Fourier transform in both variables we have

$$
\begin{equation*}
\hat{K}(\xi, \tau)=\left(2 \pi i+4 \pi^{2}|\xi|^{2}\right)^{-1} \tag{2.33}
\end{equation*}
$$

Exercise. Show that $\hat{K}$ is locally integrable near the origin, and hence defines a tempered distribution.
$\hat{K}$ is not globally integrable, however; so computing its inverse Fourier transform directly is a bit tricky. Instead, if we take the Fourier transform in the variable $x$ only,

$$
\left(\partial_{t}+4 \pi^{2}|\xi|^{2}\right) \tilde{K}(\xi, t)=\delta(t) .
$$

we can solve this $O D E$ get

$$
\tilde{K}(\xi, t)= \begin{cases}a(\xi) e^{-4 \pi^{2}|\xi|^{2} t}, & t>0 \\ b(\xi) e^{-4 \pi^{2}|\xi|^{2} t}, & t<0\end{cases}
$$

with $a(\xi)-b(\xi)=1$.
As in the previous section, there is some latitude in the choice of $a$ and $b$, but the most natural choice is the one which makes $\tilde{K}$ tempered in $t$ as well as $\xi$, namely $a=1, b=0$. So,

$$
\tilde{K}(\xi, t)=\left\{\begin{array}{l}
e^{-4 \pi^{2}|\xi|^{2} t}, t>0 \\
0 \text { otherwise }
\end{array}\right.
$$

Taking the inverse Fourier transform,

$$
K(x, t)= \begin{cases}(4 \pi t)^{-n / 2} e^{\left(|x|^{2} / 4 t\right)}, & t>0 \\ 0, & t>0\end{cases}
$$

This is a fundamental solution of the heat operator.
Remark 2.34. If we take the Fourier transform of $\tilde{K}(\xi, t)$ in the $t$ variable, we obtain

$$
\begin{aligned}
\tilde{K}(\xi, \tau) & =\int_{-\infty}^{\infty} \tilde{K}(\xi, t) e^{2 \pi \mathrm{it} \tau} d t \\
& =\int_{0}^{\infty} e^{-t\left(4 \pi^{2}|\xi|^{2}+\pi i \tau\right)} d t \\
& =\left(4 \pi^{2}|\xi|^{2}+2 \pi i \tau\right)^{-1}
\end{aligned}
$$

thus recovering formula (2.33).

Exercise. Show that this really works, i.e., the iterated Fourier transform of $K$ first in $x$, then in $t$, is the distribution Fourier transform of $K$ in all variables.

OBSERVATIONS: From the formula for $K$, we get
(1) $K(x, t)$ vanishes to infinite order as $t$ tends to 0 when $x \neq 0$ and hence is $C^{\infty}$ on $\mathbb{R}^{n+1} \backslash\{(0,0)\}$. Therefore by Theorem $2.11, \partial_{t}-\Delta$ is hypoelliptic.
(2) $K(x, t)=t^{-n / 2} K\left(\frac{x}{t}, 1\right)$. Thus if we set $K(x, 1)=\phi(x)$, then $K\left(x, \epsilon^{2}\right)$ $=\epsilon^{-n} \phi(x / \epsilon)=\phi_{\epsilon}(x)$; so, as $t$ tends to $0, K(x, t)$ tends to $\delta$, Theorem 1.6 From this, we infer
(3) If $f \epsilon L^{p}$ and if we set $u(x, t)=f *_{(x)} K(x, t)$, then

$$
\begin{aligned}
(\partial-\Delta) u & =0 \text { for } t>0 \\
u(x, 0) & =f(x) .
\end{aligned}
$$

i.e., as $t$ tends to $0,\|u(., t)-f\|_{P}$ converges to 0 . Thus we have solved the initial value problem for the homogeneous heat equation, when $f \epsilon L^{p}$. Actually, since $K(x, t)$ decreases so rapidly as $|x|$ tends to $\infty$, this works for much wider classes of $f^{\prime} s$. The convolution $f *_{(x)} K(., t)$ make sense, for example, if $|f(x)|<C e^{|x|^{2-\epsilon}}$. If $f$ is also continuous, it is not hard to see that $f *_{(x)} K(., t)$ converges to $f$ uniformly on compact sets as $t$ tends to 0 .

It is now a simple matter to solve the inhomogeneous initial value problem:

$$
\begin{aligned}
\left(\partial_{t}-\Delta\right) u & =F \text { on } \mathbb{R}^{n} \times(0, \infty), \\
u(x, 0) & =f(x) \text { on } \mathbb{R}^{n}
\end{aligned}
$$

If we take $u_{1}=F{ }_{(x, t)} K$ and

$$
u_{2}=\left(f-u_{1}(., 0)\right) *_{(x)} K
$$

then we see that

$$
\left(\partial_{t}-\Delta\right) u_{1}=F \text { on } \mathbb{R}^{n} \times(0, \infty),\left(\partial_{t}-\Delta\right) u_{2}=0
$$

on $\mathbb{R}^{n} \times(0, \infty)$, and $\left(u_{1}+u_{2}\right)(x, 0)=f(x)$. Thus $u=u_{1}+u_{2}$ solves the problem.

As another application of the fundamental solution $K$, we can derive the Weierstrass approximation theorem.

Theorem 2.35. (WEIERSTRASS) If $f$ is continuous with compact support on $\mathbb{R}^{n}$, then, for any compact $\Omega \subset \mathbb{R}^{n}$, there exists a sequence $\left(P_{m}\right)$ of polynomials such that $P_{m}$ converges to f uniformly on $\Omega$

Proof. Let $u(x, t)=\left(f *_{(x)} K(., t)\right)(x)$. Then $u(x, t)$ converges to $f(x)$ uniformly as $t$ tends to 0 . Moreover, for any $t>0$,

$$
u(x, t)=\int f(y)(4 \pi t)^{-n / 2} e^{-\sum_{j=1}^{n}\left(x_{j}-y_{j}\right)^{2} / 4 t} d y
$$

is an entire holomorphic function of $x \in \mathbb{C}^{n}$. So $u(., t)$ can be uniformly approximated on any compact set by partial sums of its Taylor series.

## 6 The Wave Operator

If we take the Fourier transform of the equation

$$
\left(\partial_{t}^{2}-\Delta\right) K=\delta(x) \delta(t)
$$

in both variables $x$ and $t$, we have, formally

$$
\begin{equation*}
\hat{K}(\xi, \tau)=\left(4 \pi^{2}|\xi|^{2}-4 \pi^{2} \tau^{2}\right)^{-1} \tag{2.36}
\end{equation*}
$$

This function $\hat{K}$ is not locally integrable, so it is not clear how to interpret it as a distribution. Again, it is better to take the Fourier transform in the $x$ variable obtaining

$$
\left(\partial_{t}^{2}+4 \pi^{2}|\xi|^{2}\right) \tilde{K}(\xi, t)=\delta(t)
$$

Solving this $O D E$,

$$
\tilde{K}(\xi, t)=a_{ \pm}(\xi) e^{2 \pi i|\xi| t}+b_{ \pm}(\xi) e^{-2 \pi i|\xi| t} \text { for } t /|t|= \pm 1
$$

and the coefficients $a_{ \pm}, b_{ \pm}$must be determined so that

$$
\tilde{K}(\xi, 0+)=\tilde{K}(\xi, 0-), \partial_{t} \tilde{K}(\xi, 0+)-\partial_{t} K(\xi, 0=)=1
$$

This gives two equations in four unknowns. In contrast to the situation with the Dirichlet problem and the heat operator, there is no way to narrow down the choices further by imposing growth restrictions on $K$ as $|t|$ tends to $\infty$. Rather, it is a characteristic feature of the wave operator that one can adapt the choice of fundamental solution to the problem at hand. The two which we shall use, called $K_{+}$and $K$ are the ones supported in the half - space $t \geq 0$ and $t \leq 0 . K_{+}$and $K_{-}$are thus determined by the requirements $a_{-}=b_{-}=0$ and $a_{+}=b_{+}=0$ respectively, from which one easily sees that

$$
\begin{aligned}
& \tilde{K}_{+}(\xi, t)=H(t) \frac{\sin 2 \pi|\xi| t}{2 \pi|\xi|} \\
& \tilde{K}_{-}(\xi, t)=-H(-t) \frac{\sin 2 \pi|\xi| t}{2 \pi|\xi|}=\tilde{K}_{+}(\xi,-t)
\end{aligned}
$$

where $H$ is the Heaviside function, i.e., the characteristic function of $[0, \infty)$. Let us compute the Fourier transforms of $\tilde{K}_{+}$and $\tilde{K}_{-}$in t to see how to make sense out of 2.36. $\tilde{K}_{+}$and $\tilde{K}_{-}$are not integrable on $t$, but it is easy to approximate them in the distribution topology by integrable functions whose Fourier transforms we can calculate. Indeed, if we set

$$
\tilde{K}_{+}^{\epsilon}(\xi, t)=e^{-2 \pi \epsilon t} H(t) \frac{\sin 2 \pi|\xi| t}{2 \pi|\xi|}, \epsilon>0,
$$

then $\tilde{K}_{+}^{\epsilon}$ is an integrable function and $\tilde{K}_{+}^{\epsilon}$ converges to $\tilde{K}_{+}$is $S^{\prime}$ as $\epsilon$ tends to 0 . Therefore $\tilde{K}_{+}=\lim _{\epsilon \rightarrow 0} \tilde{K}_{+}^{\epsilon}$ where

$$
\hat{K}_{+}(\xi, \tau)=\int_{0}^{\infty} e^{2 \pi \epsilon t-2 \pi \text { it } \tau} \frac{\sin 2 \pi|\xi| t}{2 \pi|\xi|} d t
$$

$$
\begin{aligned}
& =\int_{0}^{\infty} e^{2 \pi \epsilon t-2 \pi \text { it } \tau} \frac{\left(e^{\pi i|\xi| t}-e^{\pi i|\xi| t}\right)}{4 \pi i|\xi|} d t \\
& =\left(4 \pi^{2}\right)^{-1}\left(|\xi|^{2}-(\tau-i \epsilon)^{2}\right)^{-1}
\end{aligned}
$$

Exercise. Prove that $\hat{K}_{-}=\lim _{\epsilon \rightarrow 0} \hat{K}_{-}^{\epsilon}$ in $S^{\prime}$ where

$$
\hat{K}_{-}^{\epsilon}(\xi, \tau)=\left(4 \pi^{2}\right)^{-1}\left(|\xi|^{2}-(\tau+i \epsilon)^{2}\right)^{-1}
$$

Thus we have two distinct ways of making the function $\left[4 \pi^{2}\left(|\xi|^{2}-\right.\right.$ $\left.\left.\tau^{2}\right)\right]^{-1}$ into a tempered distribution. The difference $\tilde{K}_{+}-\tilde{K}_{-}$is of course a distribution supported on the cone $|\xi|=|\tau|$.

We now propose to use the fundamental solutions $K_{+}$and $K_{-}$to solve the Initial Value Problem or Cauchy Problem, for the operator:

$$
\begin{align*}
\left(\partial_{t}^{2}-\Delta\right) u & =f \text { on } \mathbb{R}^{n+1} \\
u(x, o) & =u_{o}(x),  \tag{2.37}\\
\partial_{t} u(x, 0) & =u_{1}(x),
\end{align*}
$$

where $u_{0}, u_{1}, f$ are given functions.
For the moment, we assume that $u_{0}, u_{1} \epsilon S$ and $f \epsilon C^{\infty}\left(\mathbb{R}_{x}^{n}\right)$ ) ( That is, $t \rightarrow f(., t)$ is a $C^{\infty}$ function with values in $S\left(\mathbb{R}^{n}\right)$ ).

Taking the Fourier transform in the variable $x$ in (2.37),

$$
\begin{aligned}
\left(\partial_{t}^{2}+4 \pi^{2}|\xi|^{2}\right) \tilde{u}(\xi, t) & =\tilde{f}(\xi, t) \\
\tilde{u}(\xi, 0) & =\tilde{u}_{0}(\xi) \\
\partial_{t} \tilde{u}(\xi, 0) & =\tilde{u}_{0}(\xi) .
\end{aligned}
$$

When $f=0$, the general solution of the $O D E$ is

$$
\begin{gathered}
\tilde{u}(\xi, t)=A(\xi) \sin 2 \pi|\xi| t+B(\xi) \cos 2 \pi|\xi| t \\
A(\xi)=\frac{\hat{u}_{1}(\xi)}{2 \pi|\xi|} \text { and } B(\xi)=\tilde{u}_{0}(\xi)
\end{gathered}
$$

when $u_{0}=u_{1}=0$, the solution is

$$
\tilde{u}(\xi, t)=\int_{0}^{t} \tilde{f}(\xi, s) \frac{\sin (2 \pi|\xi|(t-s))}{2 \pi|\xi|} d s
$$

(This may be derived by variation of parameters; in any case it is easy to check that this $u$ is in fact the solution). Therefore, the solution for the general case is given by

$$
\begin{aligned}
\tilde{u}(\xi, t) & =\hat{u}_{0}(\xi) \cos 2 \pi|\xi| t+\frac{\hat{u}_{1}(\xi)}{2 \pi|\xi|} \sin 2 \pi|\xi| t+ \\
& +\int_{0}^{t} \tilde{f}(\xi, s) \frac{\sin (2 \pi|\xi|(t-s))}{2 \pi|\xi|} d s
\end{aligned}
$$

For $t>0$, we can rewrite this as
$\tilde{u}(\xi, t)=\hat{u}_{1}(\xi) \tilde{K}_{+}(\xi, t)+\hat{u}_{0}(\xi) \partial_{t} \tilde{K}_{+}(\xi, t)+\int_{0}^{t} \tilde{f}(\xi, s) \tilde{f}(\xi, s) K_{+}(\xi, t-s) d s$.
Since $K_{+}(\xi, t-s)=0$ for $t<s$,
$\tilde{u}(\xi, t)=\hat{u}_{1}(\xi) \tilde{K}_{+}(\xi, t)+\hat{u}_{0}(\xi) \partial_{t} \tilde{K}_{+}(\xi, t)+\int_{0}^{t} \tilde{f}(\xi, s) \tilde{f}(\xi, s) K_{+}(\xi, t-s) d s$.
Take the inverse Fourier transform:

$$
\left.u(x, t)=\left(u_{1} *_{(x)} K_{+}\right)+u_{0} *_{(x)} \partial_{t} K_{+}\right)+\left(H(t) f *_{(x, t)} K_{+}\right)
$$

Likewise, for $t<0$

$$
\left.u(x, t)=-\left(u_{1} *_{(x)} K_{-}\right)-u_{0} *_{(x)} \partial_{t} K_{-}\right)+\left((H(-t) f) *_{(x, t)} K_{-}\right) .
$$

So, for arbitrary t , our solution u can be expressed as

$$
\begin{align*}
u(x, t)=\left(u_{1}\right. & *(x)\left(K_{+}-K_{)}+u_{0} *_{(x)} \partial_{t}\left(K_{+}-K_{-}\right)\right. \\
& +\left(H(t) f *_{(x, t)} K_{+}\right)+\left(H(-t) f *_{(x, t)} K_{)}\right) \tag{2.38}
\end{align*}
$$

So far, we have avoided the question of computing $K_{+}$and $K_{-}$ex- $\mathbf{5 6}$ plicitly. Indeed, since $\tilde{K}_{+}$and $\tilde{K}_{-}$are not $L^{1}$ functions, it is not an easy matter to calculate their inverse Fourier transform,

However, for the case $n=1$, we can find $K_{+}$and $K_{-}$by solving the wave wave equation directly. We have

$$
\partial_{t}^{2}-\Delta=\partial_{t}^{2}-\partial_{x}^{2}
$$

If we make the change of variables $\xi=x+t, n=x-t$, then the wave operator becomes $\partial_{t}^{2}-\partial_{x}^{2}=-4 \frac{\partial^{2}}{\partial \xi \partial \eta}$. The general solution of $\frac{\partial^{2} u}{\partial \xi \partial \eta}=0$ is given by

$$
u(\xi, n)=f(\xi)+g(\eta)
$$

where $f$ and $g$ are arbitrary functions. Therefore

$$
u(x, t)=f(x+t)+g(x-t)
$$

To solve $\left(\partial_{t}^{2}-\partial_{x}^{2}\right) u=0, u(x, 0)=u_{0}(x), \partial_{t} u(x, 0)=u_{1}(x)$ we must have

$$
\begin{aligned}
& u_{0}(x)=f(x)+g(x) \\
& u_{1}(x)=f^{\prime}(x)-g^{\prime}(x)
\end{aligned}
$$

From these equations, we have

$$
\begin{aligned}
& f^{\prime}(x)=\frac{1}{2}\left(u_{0}^{\prime}(x)+u_{1}(x)\right) \\
& g^{\prime}(x)=\frac{1}{2}\left(u_{0}^{\prime}(x)-u_{1}(x)\right)
\end{aligned}
$$

Thus $u(x, t)$ is given by

$$
u(x, t)=\frac{1}{2}\left(u_{0}(x+t)+u_{0}(x-t)\right)+\frac{1}{2} \int_{x-t}^{x+t} u_{1}(s) d s
$$

Comparing with the previous formula (2.38, we find that

$$
K_{+}(x, t)=\frac{1}{2} H(t-|x|), K_{-}(x, t)=\frac{1}{2} H(-t-|x|) .
$$

Exercise. Compute $\tilde{K}_{ \pm}$directly from these formulas.

It turns out that, for $n=2$,

$$
K_{ \pm}(x, t)=\frac{1}{2 \pi \sqrt{t^{2}-x^{2}}} H( \pm t-|x|)
$$

and, for $n=3$,

$$
K_{ \pm}(x, t)=\frac{ \pm 1}{4 \pi t} \delta( \pm t-|x|)
$$

For $n=1,2, K_{+}, K_{-}$are functions; for $n=3, K_{+}, K_{-}$are not functions but they are measures. For $n>3, K_{+}, K_{-}$are neither functions nor measures; they are more singular distributions. The exact formula for $K_{ \pm}$is rather messy and we shall not write it out; it may be read off from Theorem 5.13 and 5.14 of Folland [1] in view of our formula (2.38. The most important qualitative feature of $K_{ \pm}$, however, is that it is always supported in the light cone $\{(x, t): \pm t>|x|\}$ and we shall now prove this as a consequence of the following result.

Theorem 2.39. Suppose $u$ is a $C^{2}$ function on $\{(x, t): t \geq o\}$ such that $\left(\partial_{t}^{2}-\Delta\right) u=0$ for $t>0$ and $u=\partial_{t} u=0$ on the set $B_{0}=\left\{(x, 0):\left|x-x_{0}\right| \leq\right.$ $\left.t_{0}\right\}$. Them $u$ vanishes on $\Omega=\left\{(x, t): 0 \leq t \leq t_{0},\left|x-x_{0}\right| \leq t_{0}-t\right\}$.

Proof. Assume $u$ is real valued; otherwise, we can consider the real and imaginary parts separately. Let

$$
\begin{aligned}
& \qquad B_{t}=\left\{x:\left|x-x_{0}\right| \leq t_{0}-t\right\} \text { and } E(t)=\frac{1}{2} \int_{B_{t}}\left|\operatorname{grad}_{x, t} u\right|^{2} d x \\
& \text { i.e., } \quad E(t)=\frac{1}{2} \int_{B_{t}}\left[\left(\partial_{t} u\right)^{2}+\sum_{1}^{n}\left(\frac{\partial u}{\partial x_{j}}\right)^{2}\right] d x
\end{aligned}
$$

Then

$$
\begin{aligned}
\frac{d E}{d t}=\frac{1}{2} \int_{B_{t}} 2\left[\frac{\partial u}{\partial t} \cdot \frac{\partial^{2} u}{\partial_{t}^{2}}+\sum_{1}^{n} \frac{\partial u}{\partial x_{j}}\right. & \left.\frac{\partial^{2} u}{\partial x_{j} \partial t}\right] d x- \\
& -\frac{1}{2} \int_{\partial B_{t}}\left[\left(\frac{\partial u}{\partial t}\right)^{2}+\sum_{1}^{n}\left(\frac{\partial u}{\partial x_{j}}\right)^{2}\right] d \sigma(x) .
\end{aligned}
$$

(The second term comes from the change in the region $B_{t}$. If this is not clear, write the derivative as a limit of difference quotients and work it our).

Since

$$
\sum \frac{\partial u}{\partial x_{j}} \frac{\partial^{2} u}{\partial x_{j} \partial t}+\sum \frac{\partial^{2} u}{\partial_{x_{j}^{2}}} \frac{\partial u}{\partial t}=\operatorname{div}\left(\frac{\partial u}{\partial t} \operatorname{grad}_{x} u\right),
$$

applying the divergence theorem and using $\left(\partial_{t}^{2}-\Delta\right) u=0$, we obtain

$$
\frac{d E}{d t}=\int_{\partial B_{t}}\left[\frac{\partial u}{\partial t} \frac{\partial u}{\partial v}-\frac{1}{2}\left|\operatorname{grad}_{x, t} u\right|^{2}\right] d \sigma
$$

where $v$ is the unit normal to $B_{t}$ in $\mathbb{R}^{n}$. Now

$$
\begin{aligned}
\left|\frac{\partial u}{\partial t} \frac{\partial u}{\partial v}\right| & <\frac{1}{2}\left[\left|\frac{\partial u}{\partial t}\right|^{2}+\left|\frac{\partial u}{\partial v}\right|^{2}\right] \\
& <\frac{1}{2}\left[\left|\frac{\partial u}{\partial t}\right|^{2}+\left|\operatorname{grad}_{x} u\right|^{2}=\frac{1}{2}\left|\operatorname{grad}_{x, t} u\right|^{2}\right]
\end{aligned}
$$

Thus we see that the integrand is non-positive, and hence $\frac{d E}{d t} \leq 0$. Also $E(0)=0$ since $u=\partial_{t} u=0$ on $B_{0}$, so $E(t) \leq 0$. But $E(t) \geq 0$ by definition, so $E(t)=0$. This implies that $\operatorname{grad}_{x, t} u=0$ on $\Omega=\bigcup_{t \leq t_{0}} B_{t}$ and since $u=0$ on $B_{0}$, we conclude that $u=0$ on $\Omega$.
Corollary 2.40. Suppose $u \epsilon C^{2}$ on $\mathbb{R}^{n} \times[0, \infty),\left(\partial_{t}^{2}-\Delta\right) u=0$ for $t>$ $0, u(x, 0)=u_{0}(x), \partial_{t} u(x, 0)=u_{1}(x)$. If $\Omega_{0}=\left(\operatorname{supp} u_{0}\right) \cup\left(\operatorname{supp} u_{1}\right)$, then $\operatorname{supp} u \subset \Omega=\left\{(x, t): d\left(x, \Omega_{0}\right) \leq t\right\}$.
(The set $\Omega$ is the union of the forward light cones with vertices in $\Omega_{0}$ ).

Proof. Suppose $\left(x, t_{0}\right) \notin \Omega$. Then for some $\epsilon>0$, the set

$$
B_{0}=\left\{x: d\left(x, x_{0}\right) \leq t_{0}+\epsilon\right\} \text { is disjoint from } \Omega_{0} .
$$

Therefore, by Theorem 2.39, $u=0$ on the cone

$$
\left\{(x, t):\left|x-x_{0}\right| \leq t_{0}+\epsilon-t, 0 \leq t \leq t_{0}+\epsilon\right\} .
$$

In particular, $u=0$ on a neighbourhood of $\left(x_{0}, t_{0}\right)$, i.e., $\left(x_{0}, t_{0}\right)$ is not in the support of $u$.

## Corollary 2.41.

$$
\begin{aligned}
& \operatorname{Supp} K_{+} \subset\{(x, t): t \geq|x|\} \\
& \operatorname{Supp} K_{+} \subset\{(x, t):-t \geq|x|\}
\end{aligned}
$$

Proof. Pick $a \phi \epsilon C_{0}^{\alpha}(B(0,1))$ such that $\int \phi=1$. Put

$$
\phi_{\epsilon}(x)=\epsilon^{-n} \phi\left(\epsilon^{-1} x\right) . \text { Let } u_{\epsilon}(x, t)=\phi_{\epsilon} *_{(x)} K_{+}, t>0
$$

Then $\left(\partial_{t}^{2}-\Delta\right) u_{\epsilon}=0, u(x, 0)=0, \partial_{t} u(x, 0)=\phi_{\epsilon}(x)$ and $u_{\epsilon}$ is $C^{\alpha}$.
By the previous corollary

$$
\operatorname{supp} u_{\epsilon} \subset\{(x, t):|x| \leq t+\epsilon\}
$$

Now $u_{\epsilon}$ converges to $K_{+}$in $S^{\prime}$, as $\epsilon$ tends to 0 . Therefore,

$$
\operatorname{supp} K_{+} \subset\{(x, t):|x| \leq t\}
$$

The result for $K_{-}$follows then since $K_{-}(x, t)=K_{+}(x,-t)$.
Remarks 2.42. (i) One could also deduce the above result from our formulas for $\tilde{K}_{ \pm}$by using the Paley-Wiener theorem.
(ii) Actually for $n=3,5,7, \ldots$, it turns out that $\operatorname{supp} K_{ \pm}=\{(x, t)$ : $|x|= \pm t\}$. This is known as the Huygens principle. See Folland [1].
(iii) The distributions $K_{ \pm}$are smooth functions of $t$ (except at $t=0$ ) with values in $E^{\prime}\left(\mathbb{R}^{n}\right)$. Therefore, we now see that our formula (2.38) for the solution of the Cauchy problem makes sense even when $u_{0}, u_{1}, \epsilon D^{\prime}\left(\mathbb{R}^{n}\right)$ and $f \epsilon C\left(\mathbb{R}_{t}, D^{\prime}\left(\mathbb{R}_{x}^{n}\right)\right.$, and it is easily checked that $u$ thus defined still solves the Cauchy problem in the sense of distributions. Corollary 2.40 remains valid in this more general setting, as can seen by an approximation argument as in the proof of Corollary 2.41

EXERCISE If $\left(\partial_{t}^{2}-\Delta\right) u=f$ on $\mathbb{R}^{n+1}, u(x, 0)=u_{0}(x)$, and $\partial_{t} u(x, 0)$ $=u_{1}(x)$, figure out how $\operatorname{supp} u$ is related to $\operatorname{supp} f, \operatorname{supp} u_{0}$ and $\operatorname{supp} u_{1}$.

## Chapter 3

## $L^{2}$ Sobolev Spaces

THERE ARE MANY ways of measuring smoothness properties of functions in terms of various norms. Often it is convenient to use $L^{2}$ norms, since $L^{2}$ interacts nicely with the Fourier transforms. In this chapter, we set up a precise theory of $L^{2}$ differentiability and use it to prove Hormander's theorem on the hypoellipticity of constant coefficient differential operators.

## 1 General Theory of $L^{2}$ Sobolev Spaces

Definition 3.1. For a non-negative integer $k$, the Sobolev space $H_{k}$ is defined to be the space of all tempered distributions all of whose derivatives of order less than or equal to $k$ are in $L^{2}$.

Thus

$$
H_{k}=\left\{f \epsilon S^{\prime}: D^{\alpha} f \epsilon L^{2}\left(\mathbb{R}^{n}\right) \text { for } 0 \leq|\alpha| \leq k\right\}
$$

From the definition, we note that $f \epsilon H_{k}$ if and only if $\xi^{\alpha} \hat{f}(\xi) \epsilon L^{2}$ for $0 \leq|\alpha| \leq k$.

Proposition 3.2. $f \epsilon H_{k}$ if and only if $\left(1+|\xi|^{2}\right)^{k / 2} \hat{f} \epsilon L^{2}$.
Proof. First assume that $\left(1+|\xi|^{2}\right)^{k / 2} \hat{f} \epsilon L^{2}$. Since, for $|\xi| \geq 1$,

$$
\left|\xi^{\alpha}\right| \leq|\xi|^{k} \leq\left(1+|\xi|^{2}\right)^{k / 2} \text { for all }|\alpha| \leq k
$$

and for $|\xi|<1$,
we have
$\left|\xi^{\alpha}\right| \leq 1 \leq\left(1+|\xi|^{2}\right)^{k / 2}$ for all $|\alpha| \leq k$,
and hence

$$
\left|\xi^{\alpha} \hat{f}(\xi)\right| \leq\left(1+|\xi|^{2}\right)^{k / 2} \hat{f}(\xi)
$$

$\xi^{\alpha} \hat{f} \epsilon L^{2}$ for $|\alpha| \leq k$ which implies $f \epsilon H_{k}$.

Conversely, assume that $f \epsilon H_{k}$. Since $|\xi|^{k}$ and $\sum_{j=1}^{n}\left|\xi_{j}\right|^{k}$ are homogeneous of degree $k$ and nonvanishing for $\xi \neq 0$, we have

$$
\left(1+|\xi|^{2}\right)^{k / 2} \leq c_{0}\left(1+|\xi|^{k}\right) \leq c_{0}\left(1+c \sum_{j=1}^{n}\left|\xi_{j}\right|^{k}\right)
$$

so that

$$
\left\|\left(1+|\xi|^{2}\right)^{k / 2} \hat{f}\right\|_{2} \leq c_{0}\|f\|_{2}+c_{0} c \sum_{j=1}^{n}\left\|\partial_{j}^{k} f\right\|_{2}
$$

which shows that $\left(1+|\xi|^{2}\right)^{k / 2} \hat{f} \epsilon L^{2}$.
The characterisation of $H_{k}$ given in the above proposition immediately suggests a generalisation to non-integral values of $k$ which turns out to be very useful.
Definition 3.3. For $s \in \mathbb{R}$, we define the operator $\Lambda^{s}: S \rightarrow S$ by

$$
\left(\Lambda^{s} f\right)(\xi)=\left(1+|\xi|^{2}\right)^{s / 2} \hat{f}(\xi)
$$

64 In other words, $\Lambda^{s}=\left(I-\frac{\Delta}{4 \pi^{2}}\right)^{s / 2}$. Clearly $\Lambda^{s}$ maps $S$ continuously onto itself. We can therefore extend $\Lambda^{s}$ continuously from $S^{\prime}$ onto itself.

The Sobolev space of order $s$ is defined by

$$
H_{s}=\left\{f \epsilon S^{\prime}: \Lambda^{s} f \epsilon L^{2}\right\}
$$

We equip $H_{s}$ with the norm $\|f\|_{(s)}=\left\|\Lambda^{s} f\right\|_{2}$. If $s$ is a positive integer, the proof of Proposition 3.2 shows that $\left\|\|_{(s)}\right.$ is equivalent to the norm

$$
\|f\|=\Sigma_{0 \leq|\alpha| \leq s}\left\|D^{\alpha} f\right\|_{2}
$$

## PROPERTIES OF $H_{s}$

(i) $H_{s}$ is a Hilbert space with the scalar product defined by $(u, v)_{(s)}=$ $\left(\Lambda^{s} u, \Lambda^{s} v\right)$. Here the scalar product on the right is that in $L^{2}$. The Fourier transform is a unitary isomorphism between $H_{s}$ and the space of functions which are square integrable with respect to the measure $\left(1+|\xi|^{2}\right)^{s} d \xi$.
(ii) For every $s \in \mathbb{R}, S$ is dense in $H_{s}$.
(iii) If $s>t, H_{s} \subset H_{t}$ with continuous imbedding. In fact, for $u \epsilon H_{s}$, $\|u\|_{(t)} \leq\|u\|_{(s)}$. In particular, $H_{s} \subset L^{2}$ for $s>0$.
(iv) $D^{\alpha}$ is a bounded operator from $H_{s}$ into $H_{s-|\alpha|} S \in \mathbb{R}$.
(v) If $f \epsilon H_{-s}$, then $f$ as a linear functional on $S$ extends continuously to $H_{s}$ and $\|f\|_{(-s)}$ is the norm of $f$ in $\left(H_{s}\right)^{*}$. So we can identify $\left(H_{s}\right)^{*}$ with $H_{(-s)}$. For $f \epsilon H_{-s}, g \epsilon H_{s}$ the pairing is given by

$$
<f, g>=<\Lambda^{-s} f, \Lambda^{s} g>=\int \hat{f} \hat{g} .
$$

(If $s=0$, this identification of $H_{0}=L^{2}$ with its dual is the complex conjugate of the usual one).
(vi) The norm $\|\cdot\|_{(s)}$ is translation invariant. Indeed, if $g(x)=f\left(x-x_{0}\right)$, then $\hat{g}(\xi)=e^{2 \pi i x_{0} \cdot \xi} \hat{f}(\xi)$ and hence $\|f\|_{(s)}=\|g\|_{(s)}$.

Proposition 3.4. For $s>n / 2$, we have $\delta \epsilon H_{-s}$.
Proof. $\left(1+|\xi|^{2}\right)^{-s / 2} \hat{\delta}=\left(1+|\xi|^{2}\right)^{-s / 2} \epsilon L^{2}$, whenever $s>n / 2$. This is a consequence of the following observation :

$$
\int\left(1+|\xi|^{2}\right)^{-s} d \xi \sim 1+\int_{1}^{\infty} r^{-2 s} r^{n-1} d r<\infty \text { if and only if } s>n / 2
$$

As an immediate consequence of this proposition, we have
Corollary 3.5. For $s>n / 2+|\alpha|, D^{\alpha} \delta \epsilon H_{-s}$.

Theorem 3.6 (SOBOLEV IMBEDDING THEOREM). For $s>\frac{n}{2}+66$ $k, H_{s} \subset C^{k}$. Further, we have

$$
\begin{equation*}
\sum_{|\alpha| \leq k} \sup _{\mathbb{R}^{N}}\left|D^{\alpha} f\right| \leq C_{s k}\|f\|_{(s)} \tag{3.7}
\end{equation*}
$$

Proof. Since $S$ is dense in $H_{s}$, it suffices to prove (3.7) for $f \in S$. Let $\delta_{x}$ denote the Dirac measure at $x$. For $|\alpha| \leq k$, since $s>\frac{n}{2}+|\alpha|, D^{\alpha} \delta_{x} \epsilon H_{-s}$.

Since

$$
<D^{\alpha} \delta_{x}, f>=(-1)^{|\alpha|}<\delta_{x}, D^{\alpha} f>=(-1)^{|\alpha|}\left(D^{\alpha} f\right)(x)
$$

and $\left\|D^{\alpha} \delta_{x}\right\|_{(-s)}$ is independent of $x$,

$$
\begin{aligned}
\sum_{|\alpha| \leq k} \sup _{\mathbb{R}^{n}}\left|D^{\alpha} f(x)\right| & =\sum_{|\alpha| \leq k} \sup _{\mathbb{R}^{n}} \mid D^{\alpha} \delta_{x}, f>1 \\
& \leq \sum_{|\alpha| \leq k} \sup _{\mathbb{R}^{n}}\left|D^{\alpha} \delta_{x}\left\|_{(-s)}| | f\right\|_{(s)}=C_{s k}\|f\|_{(s)}\right.
\end{aligned}
$$

Now, given $u \epsilon H_{s}$, choose a sequence $\left(u_{j}\right)$ in $S$ such that $\left\|u-u_{j}\right\|_{(s)}$ converges to 0 , as $j$ tends to $\infty$. The above inequality with $f=u_{i}-u_{j}$ shows that ( $D^{\alpha} u_{j}$ ) is a Cauchy sequence in the uniform norm for $|\alpha| \leq k$; so its limit $D^{\alpha} u$ is continuous.

Corollary 3.8. If $u \epsilon H_{s}$ for all $s \in \mathbb{R}$, then $u \epsilon C^{\infty}$ i.e., $\bigcap_{s} H_{s} \subset C^{\infty}$.
This argument can be extended to show that if $s>\frac{n}{2}+k$, then elements of $H_{s}$ and their derivatives of order less than or equal to $k$ are not merely continuous but actually Hölder continuous.

Proposition 3.9. If $0<\alpha<1$ and $s=\frac{n}{2}+\alpha$, then $\left\|\delta_{x}-\delta_{y}\right\|_{(-s)} \leq$ $C_{\alpha}|x-y|^{\alpha}$

Proof. We have

$$
\left\|\delta_{x}-\delta_{y}\right\|_{(-s)}^{2}=\int\left|e^{-2 \pi i x \cdot \xi}-e^{-2 \pi i y \cdot \xi}\right|^{2}\left(1+|\xi|^{2}\right)^{-s} d \xi
$$

Let $R$ be a positive number, to be fixed later. When $|\xi| \leq R$ we use the estimate $\left|e^{-2 \pi i x . \xi}-e^{-2 \pi i y . \xi}\right| \leq 2 \pi|\xi \||x-y|$ (by the mean value theorem) and when $|\xi|>R$, we use $\left|e^{-2 \pi i x . \xi}-e-2 \pi i y . \xi\right| \leq 2$. Then we have

$$
\begin{aligned}
\left\|\delta_{x}-\delta_{y}\right\|_{(-s)}^{2} & \leq 4 \pi^{2}|x-y|^{2} \int_{|\xi| \leq R}|\xi|^{2}\left(1+|\xi|^{2}\right)^{-s} d \xi+4 \int_{|\xi|>R}\left(1+|\xi|^{2}\right)^{-s} d \xi \\
& \leq c\left[|x-y|^{2} \int_{0}^{R}\left(1+r^{2}\right)^{-s} r^{n+1} d r+\int_{0}^{R} r^{-2 s+n-1} d r\right] \\
& \leq c^{\prime}\left[|x-y|^{2} R^{-2 s+n+2}+R^{-2 s+n}\right] \text { as } \frac{n}{2}<s<\frac{n}{2}+1 \\
& =c^{\prime}\left[|x-y|^{2} R^{2-2 \alpha}+R^{-2 \alpha}\right]
\end{aligned}
$$

When we take $R=|x-y|^{-1}$ we get our result.
Exercise. Show that the above argument does not work when $\alpha=1$.
Instead, we get $\left\|\delta_{x}-\delta_{y}\right\|_{(-s)} \leq c\left|x-y\|\log \mid x-y\|^{1 / 2}\right.$ when $x$ is near $y$. What happens when $\alpha>1$ ?

Corollary 3.10. Let $0<\alpha<1$ and $\Lambda_{\alpha}=$ bounded functions $g$ : $\sup _{x, y} 68$ $\frac{|g(x)-g(y)|}{|x-y|^{\alpha}}<\infty$. If $s=\frac{n}{2}+\alpha+l$ and $f \epsilon H_{s}$, then $D^{\beta} f \epsilon \Lambda_{\alpha}$ for $|\beta| \leq k$.

Remark 3.11. We shall obtain an analogue of this result for $L^{p}$ norms in Chapter 5]

The following lemma will be used in several arguments hereafter.
Lemma 3.12. For all $\xi, n \in \mathbb{R}^{n}$ and $s \in \mathbb{R}$, we have

$$
\left[\frac{1+|\xi|^{2}}{1+|\eta|^{2}}\right]^{s}<2^{|s|}\left(1+|\xi-\eta|^{2}\right)^{|s|}
$$

Proof. $|\xi| \leq|\eta|+|\xi-\eta|$ gives

$$
\begin{gathered}
|\xi|^{2}<2\left(|\eta|^{2}+|\xi-\eta|^{2}\right) \text { so that } \\
\left(1+|\xi|^{2}\right)<2\left(1+|\eta|^{2}\right)\left(1+|\xi-\eta|^{2}\right)
\end{gathered}
$$

If $s>0$, then raise both sides to the $s^{\text {th }}$ power. If $s<0$, interchange $\xi$ and $\eta$ and raise to the $-s^{\text {th }}$ power.

Proposition 3.13. If $\phi \epsilon S$, then the operator $f \rightarrow \phi f$ is bounded for all $s$.

Proof. The operator $f \rightarrow \phi f$ is bounded on $H_{s}$ if and only if the operator of $g \rightarrow \Lambda^{s} \phi \Lambda^{-s} g$ is bounded on $L^{2}$, as one sees by
$69 \quad$ setting $g=\wedge^{s} f$. But

$$
\begin{aligned}
\left(\wedge^{s} \phi^{-s} g\right)^{\wedge}(\xi) & =\left(1+|\xi|^{2}\right)^{s / 2}\left(\phi \wedge^{-s} g\right)^{\wedge}(\xi) \\
& =\left(1+|\xi|^{2}\right)^{s / 2}\left[\hat{\phi}(\xi) *\left(\wedge^{-s} g\right)(\xi)\right] \\
& =\left(1+|\xi|^{2}\right)^{s / 2} \int \hat{\phi}(\xi-\eta)\left(1+|\eta|^{2}\right) \hat{g}(\eta) d \eta \\
& =\int \hat{g}(\eta) K(\xi, \eta) d \eta
\end{aligned}
$$

where

$$
K(\xi, \eta)=\left(1+\left.\eta\right|^{2}\right)^{-s / 2}\left(1+|\xi|^{2}\right)^{s / 2} \hat{\phi}(\xi-\eta)
$$

By lemma 3.12

$$
|K(\xi, \eta)|<2^{|s| / 2}\left(1+|\xi-\eta|^{2}\right)^{|s| / 2}|\hat{\phi}(\xi-\eta)|
$$

Therefore, since $\hat{\phi}$ is rapidly decreasing at $\infty$,

$$
\begin{aligned}
& \int|K(\xi, \eta)| d \xi \leq c \text { for every } \eta \in \mathbb{R}^{n} \\
& \int|K(\xi, \eta)| d \eta \leq c \text { for every } \eta \in \mathbb{R}^{n}
\end{aligned}
$$

Thus from Theorem 1.1 the operator with kernel $K$ is bounded on $L^{2}$. Hence our proposition is proved.

The spaces $H_{s}$ are defined on $\mathbb{R}^{n}$ globally by means of the Fourier transform. Frequently, it is more appropriate to consider the following versions of these spaces.

Definition 3.14. If $\Omega \subset \mathbb{R}^{n}$ is open and $s \in \mathbb{R}$, we define $H_{s}^{\mathrm{loc}} \Omega=\left\{f \in \mathcal{D}^{\prime}\right.$ $(\Omega): \forall \Omega^{\prime} \Subset \Omega, \exists g_{\Omega}, \epsilon H_{s}$ such that

$$
\left.g_{\Omega^{\prime}}=f \text { on } \Omega^{\prime}\right\}
$$

Proposition 3.15. $f \epsilon H_{s}^{\text {loc }}(\Omega)$ if and only if $\phi f \epsilon H_{s}$ for $\phi \epsilon C_{0}^{\infty}(\Omega)$.
Proof. If $f \epsilon H_{s}^{\mathrm{loc}}(\Omega)$ and $\phi \epsilon C_{0}^{\infty}(\Omega)$, then there exists $g \epsilon H_{s}$ such that $f=$ $g$ on supp $\phi$. Therefore $\phi f=\phi g \epsilon H_{s}$ by proposition 3.13.

Conversely, if $\phi f \epsilon H_{s}$ for all $\phi \epsilon C_{0}^{\infty}\left(\Omega^{\prime}\right)$, and $\Omega^{\prime} \Subset \Omega$, choose $\phi \epsilon C_{0}^{\infty}(\infty)$ with $\phi \equiv 1$ on $\Omega^{\prime}$. Then $\phi f \epsilon H_{s}$ and $f=\phi f$ on $\Omega^{\prime}$.
Corollary 3.16. If $L=\sum_{|\propto| \leq k} a_{\propto}(x) D^{\propto}$ with $a_{\alpha} \epsilon C^{\infty}(\Omega)$, then $L$ maps $H_{s}^{\mathrm{loc}}(\Omega)$ into $H_{s-k}^{\mathrm{loc}}(\Omega)$ for all $s \in \mathbb{R}$.

It is a consequence of the Arzela-Ascoli theorem that if $\left(u_{j}\right)$ is a sequence of $C^{k}$ functions such that $\left|u_{j}\right|$ and $\left|\partial^{\alpha} u_{j}\right|(|\alpha| \leq k)$ are bounded on compact set uniformly in $j$, there exists a subsequence $\left(v_{j}\right)$ of $\left(u_{j}\right)$ such that ( $\partial^{\alpha} v_{j}$ ) converges uniformly on compact set for $|\alpha| \leq k-1$. In particular, if the $u_{j}^{\prime} \mathrm{s}$ are supported in a common compact set, then $\left(\partial^{\alpha} v_{j}\right)$ converges uniformly.

There is an analogue of this result for $H_{s}$ spaces.
Lemma 3.17. Suppose $\left(u_{k}\right)$ is a sequence of $C^{\infty}$ functions supported in a fixed compact set $\Omega$ such that $\sup _{k}\left\|u_{k}\right\|_{(s)}<\infty$. Then there exists a subsequence which converges in the $H_{t}$ norm for all $t<s$.
Proof. Pick a $\phi \epsilon C_{0}^{\infty}$ such that $\phi=1$ on $\Omega$ so that $u_{k}=\phi u_{k}$ and hence $\mathbf{7 1}$ $\hat{u}_{k}=\hat{\phi}_{*} \hat{u}_{k}$. Then

$$
\begin{aligned}
\left(1+|\xi|^{2}\right)^{s / 2}\left|\hat{u}_{k}(\xi)\right| & =\left(1+|\xi|^{2}\right)^{s / 2}\left|\int \hat{\phi}(\xi-\eta) \hat{u}_{k}(\eta) d \eta\right| \\
& \leq \int\left|\hat{\phi}(\xi-\eta) \| \hat{u}_{k}(\eta)\right| 2^{|s| / 2}\left(1+|\eta|^{2}\right)^{s / 2}\left(1+|\xi-\eta|^{2}\right)^{\mid s / 2} d \eta \\
& \leq 2^{|s| / 2}\|\phi\|_{(|s|)}\left\|u_{k}\right\|_{(s)} \leq c-{ }_{1} \text { independent of } k
\end{aligned}
$$

Likewise, we have

$$
\left(1+|\xi|^{2}\right)^{s / 2}\left|\partial_{j} \hat{u}_{k}(\xi)\right| \leq 2^{|s| / 2}\left\|2 \pi x_{j} \phi(x)\right\|_{|s|}\left\|u_{k}\right\|_{(s)} \leq c_{2}
$$

independently of $k$. Therefore, by the Arzela-Ascoli theorem there exists a subsequence $\left(\hat{v}_{k}\right)$ of $\left(\hat{u}_{k}\right)$ which converges uniformly on compact sets. For $t \leq s$,

$$
\begin{aligned}
\left\|v_{j}-v_{k}\right\|_{(t)}^{2}= & \int\left(1+|\xi|^{2}\right)^{t}\left|\hat{v}_{j}-\hat{v}_{k}\right|^{2} d \xi \\
= & \int_{|\xi| \leq R}\left(1+|\xi|^{2}\right)^{t}\left|\hat{v}_{j}-\hat{v}_{k}\right|^{2} d \xi+\int_{|\xi|>R}\left(1+|\xi|^{2}\right)^{t}\left|\hat{v}_{j}-\hat{v}_{k}\right|^{2} d \xi \\
< & \left(1+R^{2}\right)^{\max (t, 0)} \sup _{|\xi| \leq R}\left|\hat{v}_{j}(\xi)-\hat{v}_{k}(\xi)\right|^{2} \int_{|\xi| \leq R} d \xi+ \\
& +\left(1+R^{2}\right)^{t-s} \int_{|\xi|>R}\left[\left(1+|\xi|^{2}\right)^{s}\left|\hat{v}_{j}(\xi)-\hat{v}_{k}(\xi)\right|^{2} d \xi\right] \\
\leq & c\left(l+R^{2}\right)^{n+|t|} \sup _{|\xi| \leq R}\left|\hat{v}_{j}(\xi)-\hat{v}_{j}(\xi)\right|^{2}+\left(1+R^{2}\right)^{t-s}\left\|v_{j}-v_{k}\right\|_{(s)}^{2}
\end{aligned}
$$

Given $\epsilon>0$, choose $R$ large enough so that the second term is less than $\epsilon / 2$ for all $j$ and $k$. This is possible since $\left\|v_{j}-v_{k}\right\|_{(s)} \leq c$ and $t-s<0$. Then for $j$ and $k$ large enough the first term is less than $\epsilon / 2$, since $\left(\hat{v}_{k}\right)$ converges uniformly on compact sets. Thus we see that $\left(v_{k}\right)$ is a Cauchy sequence in $H_{t^{\prime}}$ and since $H_{t}$ is complete we are done.

Remark 3.12. Lemma 3.17 is false, if we do not assume that all the $u_{k}^{\prime} s$ have support in a fixed compact set. For example, for $u \in C_{0}^{\infty}$ and $x_{k} \in \mathbb{R}^{n}$ with $\left|x_{k}\right|$ tending to $\infty$, define $u_{k}(x)=u\left(x-x_{k}\right)$. Then the invariance of $H_{s}$ norms shows that $\left\|u_{k}\right\|_{(s)}=\|u\|_{(s)}$. But no subsequence of $\left(u_{k}\right)$ converges, in any $H_{t}$. For, if a subsequence $\left(v_{k}\right)$ of $\left(u_{k}\right)$ converges, it must converge to 0 , since $u_{k}$ converges to 0 in $S^{\prime}$. But then $\lim \left\|v_{k}\right\|_{(t)}=0$ which is not the case.

Theorem 3.20 (RELLICH THEOREM). Let $H_{s}^{0}(\Omega)$ be the closure of $C_{0}^{\infty}(\Omega)$ in $H_{s}$. If $\Omega$ is bounded and $t<s$, the inclusion $H_{s}^{0}(\Omega) \hookrightarrow H_{t}$ is compact, i.e., bounded sets in $H_{s}^{0}(\Omega)$ are relatively compact in $H_{t}$.

Proof. Let $\left(u_{k}\right)$ be a sequence in $H_{s}^{0}(\Omega)$. To each $k$, find a $v_{k} \epsilon C_{0}^{\infty}(\Omega)$
such that $\left\|u_{k}-v_{k}\right\|_{(s)} \leq \frac{1}{k}$. Then we have $\left\|v_{k}\right\|_{(s)} \leq\left\|u_{k}\right\|_{(s)}+\mid \leq c$ (independent of $k$.) Therefore, by lemma 3.17 a subsequence $\left(w_{k}\right)$ of $\left(v_{k}\right)$ exists such that $\left(w_{k}\right)$ converges in $H_{t}$. If $\left(u_{k}^{\prime}\right)$ is the subsequence of $\left(u_{k}\right)$ corresponding to the sequence $\left(w_{k}\right)$, we have

$$
\begin{aligned}
\left\|u_{i}^{\prime}-u_{j}^{\prime}\right\|_{(t)} & <\left\|u_{i}^{\prime}-w_{i}\right\|_{(t)}+\left\|w_{i}-w_{j}\right\|_{(t)}+\left\|w_{j}-u_{j}^{\prime}\right\|_{(t)} \\
& <\frac{1}{i}+\frac{1}{j}+\left\|w_{i}-w_{j}\right\|_{(t)} \rightarrow 0 \text { as } \| i, j \rightarrow \infty .
\end{aligned}
$$

Hence $\left(u_{k}^{\prime}\right)$ converges in $H_{t}$.
I the proof of the next theorem, we will use the technique of complex interpolation, which is based on the following result from elementary complex analysis called 'Three lines lemma'.
Lemma 3.20. Suppose $F(z)$ is analytic in $o<\operatorname{Re} Z<1$, continuous and bounded on $0 \leq \operatorname{Re} Z \leq 1$. If $|F(1+i t)| \leq c_{0}$ and $|F(l+i t)| \leq c_{1}$, then $F(s+i t) \leq c_{0}^{1-s} c_{1}^{s}$, for $0<s<1$.
Proof. If $\epsilon>0$, the function

$$
g_{\epsilon}(z)=c_{0}^{z-1} c_{1}^{z-1} e^{\epsilon\left(z^{2}-z\right)} f(z)
$$

satisfies the hypotheses with $c_{0}$ and $c_{1}$ replaced by 1 , and also $\left|g_{\epsilon}(z)\right|$ converges to 0 as $|I m z| \rightarrow \infty$ for $0 \leq \operatorname{Re} Z \leq 1$. From the maximum modulus principle, it follows that $|g(z)| \leq 1$ for $0 \leq \operatorname{Re} z \leq 1$ and letting $\epsilon$ tend to 0 , we obtain the desired result.

Theorem 3.21. Suppose that $-\infty<s_{0}<s_{1}<\infty$ and $T$ is a bounded linear operator $H_{s_{0}}$ such that $T \mid H_{s_{1}}$ is bounded on $H_{s_{1}}$. Then $T \mid H_{s}$ is bounded on $H_{s}$ for all $s$ with $s_{0} \leq s \leq s_{1}$.

Proof. Our hypothesis means that

$$
\wedge^{s_{0}} T \wedge^{-s_{0}} \text { and } \wedge^{s_{1}} T \wedge^{-s_{1}}
$$

are bounded operators on $L^{2}$. For $0 \leq \operatorname{Rez} \leq 1$ we define

$$
s_{z}=(l-z) s_{0}+z s_{1} \text { and } T_{z}=\wedge^{s_{z}} T \wedge^{-s_{z}}
$$

Then what we wish to prove is that $T_{z}$ is bounded on $L^{2}$ for $0 \leq z \leq$ 1. Observe that when $w=x+i y, \wedge^{w}, \wedge^{x} \wedge^{i y}$ and

$$
\begin{aligned}
\left(\wedge^{i y} f\right)^{\wedge}(\xi) & =\left(1+|\xi|^{2}\right)^{i y / 2} \hat{f}^{2}(\xi) \text { so that } \\
\mid\left(\wedge^{i y} f\right)^{\wedge}(\xi) & =|\hat{f}(\xi)| .
\end{aligned}
$$

Thus $\wedge^{i y}$ is unitary on $H^{s}$ for all $s$.
For $\phi, \psi \in S$, we define

$$
F(z)=\int\left(T_{z} \phi\right) \psi=<\wedge^{s_{z}} T \wedge^{-s_{z}} \phi, \psi>
$$

Then

$$
\begin{aligned}
|F(z)| & =\left|<\wedge^{s_{z}} T \wedge^{-s_{z}} \phi, \psi>\right| \\
& =\mid<T \wedge^{s_{z}} \phi, \wedge^{-s_{z}} \psi>1 \\
& \leq\left\|T \wedge^{-s_{z}} \phi\right\|_{\left(s_{0}\right)}\left\|\wedge^{-s_{z}} \psi\right\|_{\left(-s_{0}\right)} \\
& \leq c\left\|\wedge^{-s_{z}} \phi\right\|_{s_{0}}\left\|\wedge^{s_{z}}\right\|_{\left(-s_{0}\right)} \\
& \leq c\|\phi\|_{\left(s_{0}-s_{1}\right) R e}\|\psi\|_{\left(s_{1}-s_{0}\right) R e z} \\
& \leq c\|\phi\|_{(0)}\|\psi\|_{s_{1}-s_{0}}
\end{aligned}
$$

$F(z)$ is clearly an analytic function of $z$ for $0<\operatorname{Re} z<1$. Further, by our hypothesis on $T$, when $\operatorname{Rez}=0$, we have

$$
|F(z)| \leq c_{0}\|\phi\|_{(0)}\|\psi\|_{(0)}
$$

and when $\operatorname{Rez}=1$, we have

$$
|F(z)| \leq c_{1}\|\phi\|_{(0)}\|\psi\|_{(0)} .
$$

Therefore, by the Three lines lemma

$$
|F(z)| \leq c_{0}^{1-z} c_{0}^{z}\|\phi\|_{(0)}\|\psi\|_{(0)} \text { for } 0<z<1 .
$$

Finally, by the self duality of $H_{0}=L^{2}$, this gives

$$
\left\|T_{z} \phi\right\| \leq c_{0}^{1-z} c_{0}^{z}\|\phi\|_{(0)}
$$

which completes the proof.

Remark 3.22. The same proof also yields the following more general result:

Suppose $\infty<s_{0}<s_{1}<\infty<t_{0}<t_{1}<\infty$. If $T$ is a bounded linear operator from $H_{s_{0}}$ to $H_{t_{0}}$ whose restriction to $H_{s_{1}}$ to $H_{t_{1}}$, then the restriction of $T$ to $H_{t_{\theta}}$ is bounded from $H_{t_{\theta}}$ to $H_{t_{\theta}}$ for $0<\theta<1$ where $s_{\theta}=(1-\theta) s_{0}+\theta s_{1}$ and $t_{\theta}=(1-\theta) t_{0}+\theta t_{1}$.

As a consequence of this result, we obtain an easy proof that $H_{s}^{\text {loc }}$ is invariant under smooth coordinate changes.

Theorem 3.23. Suppose $\Omega$ and $\Omega^{\prime}$ are open subsets of $\mathbb{R}^{n}$ and $\phi: \Omega \rightarrow$ $\Omega^{\prime}$ is a $C^{\infty}$ diffeomorphism. Then the mapping $f \rightarrow$ foф maps $H_{s}^{\mathrm{loc}}\left(\Omega^{\prime}\right)$. continuously onto $H_{s}^{\mathrm{loc}}(\Omega)$.

Proof. The statement of the theorem is equivalent to the assertion that for any $\phi \epsilon C_{0}^{\infty}\left(\Omega^{\prime}\right)$, the map $T f=(\phi f) \circ \phi$ is bounded on $H_{s}$ for $s \in \mathbb{R}$. If $s=0,1,2, \ldots$, this follows from the chain rule and the fact that $H_{s}=$ $\left\{f: D^{\propto} f \epsilon L^{2}\right.$ for $\left.|\propto| \leq s\right\}$. By Theorem 3.21 it is true for all $s \geq 0$. But the adjoint of $T$ is another map of the same form :
$T^{*} g=(\psi g)_{\circ} \psi$ where $\psi=\theta^{-1}$ and $\psi=\phi|J| \circ \Phi, J$ being the Jacobian determinant of $\Phi^{-1}$. Hence $T^{*}$ is bounded on $H_{s}$ for all $s \geq 0$ and by duality of $H_{s}$ and $H_{-s}$, this yields the boundedness of $T$ on $H_{s}$ for $s<0$.

Finally, we ask to what extend the $H_{s}$ spaces include all distributions. Globally they do not, since, if $f \epsilon H_{s}$, then $f$ is tempered and $\hat{f}$ is a function. But locally they do, as we see from the following result.

Proposition 3.24. Every distribution with compact support lies in some $H_{s}:$ i.e., $E^{\prime} \subset \bigcup_{s \in \mathbb{R}} H_{s}$.

Proof. If $f \epsilon E^{\prime}$, then it is a continuous linear functional on $C^{\infty}$. Therefore, there exists a constant $c>0$, a compact set $K$, and a nonnegative integer $k$ such that

$$
\left|<f, \phi>\left|\leq c \sum_{|\propto| \leq k} \sup _{k}\right| D^{\infty} \phi\right| \text { for all } \phi \epsilon C^{\infty}
$$

i.e., $\left|<f, \phi>\left|\leq c \sum_{|\alpha| \leq \mid k} \sup _{\mathbb{R}^{n}}\right| D^{\alpha} \phi\right|$ for all $\phi \epsilon C^{\infty}$.

By the Sobolev imbedding theorem,

$$
\sum_{|\alpha| \leq k} \sup _{\mathbb{R}^{n}}\left|D^{\alpha} \phi\right| \leq c^{\prime}\|\phi\|_{\left(k+\frac{n}{2}+\epsilon\right)} \text { for } \epsilon>0
$$

Therefore, $|<f, \phi>| \leq c^{\prime \prime} \|_{\left(k+\frac{n}{2}+\epsilon\right)}$ for all $\phi \epsilon S$ since $S$ is dense in $H_{\left(k+\frac{n}{2}+\epsilon\right)}$, this shows that $f$ is a continuous linear functional on $H_{\left(k+\frac{n}{2}+\epsilon\right)}$. Hence $f \epsilon H_{-\frac{n}{2}-k-\epsilon}$.

Corollary 3.25. If $f \in D^{\prime}(\Omega)$ and $\Omega^{\prime}$ has compact closure in $\Omega$ then there exists $s$ in $\mathbb{R}$ such that $f \in H_{s}^{\mathrm{loc}}\left(\Omega^{\prime}\right)$.

## 2 Hypoelliptic Operators With Constant Coefficients

We now apply the machinery of Sobolev spaces to derive a criterion for the hypoellipticity of constant coefficient differential operators. First, we have a few preliminaries.

Definition 3.26. Let $P$ be a polynomial in $n$ variables. For a multi-index $\alpha, P^{(\alpha)}$ will ne defined by $P^{(\alpha)}(\xi)=\left(\frac{\partial}{\partial \xi}\right)^{\alpha} P(\xi)$.
Proposition 3.27. LEIBNIZ RULE When $f \epsilon C^{\infty}, g \epsilon D^{\prime}$ and $P(D)$ is a constant-coefficient partial differential operator of order $k$, we have

$$
\left.P(D)(f g)=\sum_{|\alpha| \leq k} \frac{1}{\alpha!}\left(P^{(\alpha)}(D) g\right) D^{\alpha} f\right) .
$$

The proof of this proposition is left as an exercise to the reader.
78 Definition 3.28. We say that a polynomial $P$ satisfies condition $(H)$ if there exists a $\delta>0$ such that

$$
\frac{\left|P^{(\alpha)}(\xi)\right|}{|P(\xi)|}=0\left(|\xi|^{-\delta|\alpha|}\right) a s|\delta| \rightarrow \infty, \forall \propto .
$$

Theorem 3.29 (HÖRMANDER). If $P$ satisfies condition $(H)$, then $P(D)$ is hypoelliptic. More precisely, if $f$ is in $\left.D^{\prime} \Omega\right)$ and $P(D) f H_{s+k \delta}^{\mathrm{loc}}(\Omega)$, where $\delta$ is as in condition $(H)$ and $k$ is the degree of $P$.

Proof. We first observe that the second assertion implies the first since $C^{\infty}(\Omega)=\bigcap_{s \in \mathbb{R}} H_{s}^{\text {loc }}(\Omega)$ by Corollary 3.8

Suppose therefore that $P(D) f \epsilon H_{s}^{\mathrm{loc}}(\Omega), f \epsilon D^{\prime}(\Omega)$. Given $\phi \epsilon C_{\circ}^{\infty}(\Omega)$ we have to prove that $\phi f \epsilon H_{s+k \delta}$. Let $\Omega^{\prime}$, be an open set such that $\Omega^{\prime} \Subset \Omega$ and $\operatorname{supp} \phi \subset \Omega^{\prime}$. By Corollary 3.25 therefore exists $t$ in $\mathbb{R}$ such that $f \epsilon H_{t}^{\text {loc }}\left(\Omega^{\prime}\right)$. By decreasing $t$, we can some that $t=s+k-$ $1-m \delta$ for assume positive integer $m$. Set $\phi_{m}=\phi$ and then choose $\phi_{m-1}, \phi_{m-2}, \ldots, \phi_{0}, \phi-1$ in $C_{0}^{\infty}\left(\Omega^{\prime}\right)$ such that $\phi_{j}=1$ on $\operatorname{supp} \phi_{j+1}$.

Then $\phi_{j} P(D) f \epsilon H_{S} \subset H_{t-k+l+j \delta}$ for $0 \leq j \leq m$ and $\phi_{-1} f \epsilon H_{t}$. Now

$$
P(D)\left(\phi_{0} f\right)=\phi_{0} P(D) f+\sum_{\alpha \neq 0} \frac{1}{\propto!} P^{\propto}(D)\left(\phi_{-1} f\right) D^{\alpha_{\phi_{0}}}
$$

since $\theta_{-1}=1$ on the support of $\phi_{0}$. So, $P(D)\left(\phi_{0} f\right) \epsilon H_{r-k+1}$. This means that

$$
\int\left(1+|\xi|^{2}\right)^{t-k+1}\left|P(\xi)\left(\phi_{0} f\right)^{\wedge}(\xi)\right|^{2} d \xi<\infty
$$

By condition (H)

$$
\left.\int\left(1+|\xi|^{2}\right)^{t-k+l+\delta|\delta|} \mid P^{\alpha \mid}(\xi)\left(\phi_{0} f\right)^{( } \xi\right)\left.\right|^{2} d \xi<\infty
$$

This implies that $P^{(\alpha)}(D)\left(\phi_{\circ} f\right) \epsilon H_{t}-k+l+\delta|\delta|$.
Next,

$$
P(D)\left(\phi_{1} f\right)=\phi_{1} P(D) f+\sum_{\alpha \neq 0} \frac{1}{\propto!} p^{(\alpha)}(D)\left(\phi_{o} f\right) D^{\propto}\left(\phi_{1} f\right)
$$

since $\phi_{0}=1$ on the support of $\phi_{1}$, so $P(D)\left(\phi_{1} f\right) \epsilon H_{t-k+1+\delta}$. By the same argument same argument as above,

$$
P^{(D)}\left(\phi_{1} f\right) \epsilon H_{t+k+l+\delta(1+|\propto|)} .
$$

Continuing inductively, we obtain $P(D)\left(\phi_{j} f\right) \epsilon H_{t-k+l+j \delta}$, which implies that

$$
P^{(\propto)}(D)\left(\phi_{j} f\right) \epsilon H_{t-k+l+\delta(j+|\propto|)} .
$$

For $j=m$, we above

$$
P^{(\alpha)}(D)\left(\phi_{m} f\right) \epsilon H_{t-k+l+\delta(m+(\delta))}=H_{s+\delta|\propto|} .
$$

If

$$
P(\xi)=\sum_{|\propto| \leq k} a_{\propto} \xi^{\propto}
$$

choose $\propto$ with $|\propto|=k$ such that $a_{\propto} \neq 0$. Then $P^{\propto}(\xi)=\propto!a_{\propto} \neq 0$. whence $\phi f=\phi_{m} f \epsilon H_{s+k \delta}$ and we are done.
$\mathbf{8 0}$ Remark 3.30. The condition $(\mathrm{H})$ is equivalent to the following apparently weaker condition

$$
\left(H^{\prime}\right): \frac{\left|p^{(\alpha)(\xi)}\right|}{|P(\xi)|} \rightarrow 0 \text { as }|\xi| \rightarrow \infty \text { for } \propto \neq 0
$$

Condition $\left(H^{\prime}\right)$ is in turn equivalent to

$$
\left(H^{\prime \prime}\right):|\operatorname{Im} \zeta| \rightarrow \infty, \text { in the set }\left\{\zeta \in \mathbb{C}^{n}: P(\zeta)=0\right\}
$$

The converse of Hörmander's theorem is also true, i.e., hypoellipticity implies condition (H).

The proofs of these assertions can be found in Hörmander [6]. The logical order of the proofs is

$$
(H) \Rightarrow \text { hypoellipticity } \Rightarrow\left(H^{\prime \prime}\right) \Rightarrow\left(H^{\prime}\right) \Rightarrow(H)
$$

The implication $\left(H^{\prime}\right) \Rightarrow(H)$ requires the use of some results from (semi) algebraic geometry.

Definition 3.31. $P(\xi)=\sum_{|\propto| \leq k} a_{\propto} \xi^{\propto}$ is called elliptic if $\sum_{|\propto|=k} a_{\propto} \xi^{\propto} \neq 0$ for every $\xi \neq 0$.

## EXERCISES

1. Prove that $P$ is elliptic if and only if $|P(\xi)| \geq c_{\circ}|\xi|^{k}$ for large $|\xi|$.
2. Prove that $P$ is elliptic if and only if $P$ satisfies condition $(H)$ with $\delta=1$.
3. Prove that no $P$ satisfies condition $(H)$ with $\delta>1$. (Hint: If $|\propto|=k, P^{(\alpha)}$ is a constant) .
4. Let $P$ be elliptic and real valued. Define $Q$ on $\mathbb{R}^{n+1}$ by $Q(\xi, \tau)=$ $2 \pi i \tau+P(\xi)$, so that $Q\left(D_{x}, D_{t}\right)=\partial_{t}+P\left(D_{x}\right)$. Show that $Q$ satisfies condition $(H)$ with $\delta=1 / k$ where $k$ is the degree of $P$ and that $1 / k$ is the best possible value of $\delta$.
(Hint : Consider the regions $|\xi|^{k} \leq|\tau|$ and $|\tau| \leq|\xi|^{k}$ separately).

## Chapter 4

## Basic Theory of Pseudo Differential Operators

## 1 Representation of Pseudo differential Operators

Let $L=\sum_{|\alpha| \leq k} a_{\alpha}(x) D^{\alpha}$ be a partial differential operator with $C^{\infty}$ coeffi- 82 cients on $\Omega$. Using the Fourier transform, we can write

$$
\begin{aligned}
(L u)(x) & =\sum_{|\alpha| \leq k} a_{\alpha}(x) \int e^{2 \pi i x \cdot \xi} \hat{u}(\xi) \xi^{\alpha} d \xi \\
& =\int e^{2 \pi i x . \xi} p(x, \xi) \hat{u}(\xi) d \xi
\end{aligned}
$$

where $p(x, \xi)=\sum_{|\alpha| \leq k} a_{\alpha}(x) \xi^{\alpha}$. This representation suggests that $p(x, \xi)$ can be replaced by more general functions. So, we make the following definition.

Definition 4.1. For an open set $\Omega \subset \mathbb{R}^{n}$ and a real number $m$ we define $S^{m}(\Omega)$, the class of symbols of order $m$ on $\Omega$, by

$$
S^{m}(\Omega)=\left\{p \epsilon C^{\infty}\left(\Omega \times \mathbb{R}^{n}\right): \forall \alpha, \beta, V \Omega^{\prime} \subset \Omega, c=c_{\alpha \beta \Omega^{\prime}}\right.
$$

such that $\left.\sup _{x \in \Omega^{\prime}}\left|D_{x}^{\beta} D_{\xi}^{\alpha} p(x, \xi)\right| \leq c(1+|\xi|)^{m-|\alpha|}\right\}$.

We note that $S^{m}(\Omega) \subset S^{m}$ whenever $m<m^{\prime}$, and we set $S^{-\infty}(\Omega)=$ $\bigcap_{m \in \mathbb{R}}$ $S^{m}(\Omega)$.

## Examples

(i) Let $p(x, \xi)=\sum_{|\alpha| \leq k} a_{\alpha}(x) \xi^{\alpha}$ with $a_{\alpha} \epsilon C^{\infty}(\Omega)$. Then $p \epsilon S^{k}(\Omega)$.
(ii) Let $p(x, \xi)=\sum_{j=1}^{N} a_{j}(x) f_{j}(\xi)$ where $a_{f} \epsilon C^{\infty}(\Omega)$ and $f_{j} \epsilon C^{\infty}\left(\mathbb{R}^{n}\right)$ is homogeneous of degree $m_{j}$ large $\xi$ : that is,

$$
f_{i}(r \xi)=r^{m_{j}} f_{j}(\xi) \text { for }|\xi| \geq c, r \geq 1
$$

In this case, $p \in S^{m}(\Omega)$, where $m=\max _{1 \leq j \leq N}\left\{m_{j}\right\}$.
(iii) Let $p(x, \xi)=\left(1+|\xi|^{2}\right)^{s / 2}$. This $p$ belongs to $S^{s}\left(\mathbb{R}^{n}\right)$.
(iv) Let $p(x, \xi)=\phi(\xi) \sin \log |\xi|$ with $\phi \epsilon C^{\infty}, \phi=0$ near the origin and $\phi=1$ when $|\xi| \geq 1$. Then $p \in S^{o}\left(\mathbb{R}^{n}\right)$.

Remark 4.2. We observe that when $p \epsilon S^{m}(\Omega), D_{x}^{\beta} D_{\xi}^{\alpha} p \epsilon S^{m-|\alpha|}(\Omega)$.
Further, if $p \epsilon S^{m_{1}}(\Omega)$ and $q \epsilon S^{m_{2}}(\Omega)$, then $p+q \epsilon S^{m}(\Omega)$, where $m=$ $\max \left\{m_{1}, m_{2}\right\}$ and $p q \epsilon S^{m_{1}+m_{2}}(\Omega)$.

Remark 4.3. Our symbol classes $S^{m}(\Omega)$ are special cases of Hormander's classes $S_{\rho, \delta}^{m}(\Omega)$. Namely, for $0 \leq \delta \leq \rho \leq 1$, and $m \in \mathbb{R}$,

$$
S_{\rho, \delta}^{m}(\Omega)=\left\{p \epsilon C^{\infty}\left(\Omega \times \mathbb{R}^{n}\right): \forall \alpha, \beta, V \Omega^{\prime} \Subset \Omega, c=c_{\alpha \beta \Omega^{\prime}}\right.
$$

such that

$$
\left.\sup _{x \in \Omega^{\prime}}\left|D_{x}^{\beta} D_{\xi}^{\alpha} p(x, \xi)\right| \leq c(1+|\xi|)^{m-\rho|\alpha|+\delta|\beta|}\right\} .
$$

In this terminology, $S^{m}(\Omega)=S_{1,0}^{m}(\Omega)$.
Definition 4.4. For $p \epsilon S^{m}(\Omega)$, we define the operator $p(x . D)$ on the do$\operatorname{main} C_{o}^{\infty}(\Omega)$ by

$$
p(x, D) u(x)=\int e^{2 \pi i x . \xi} p(x, \xi) \hat{u}(\xi) d \xi, u \epsilon C_{o}^{\infty}(\Omega)
$$

(Sometimes, we shall denote $p(x, D)$ by $p$ ). Operators of the form $p(x, D)$ with $p \in S^{m}(\Omega)$ are called pseudo differential operators of order $m$ on $\Omega$.

The set of all pseudo differential operators of order $m$ on $\Omega$ will denoted by $\Psi^{m}(\Omega)$. For brevity, we will sometimes write " $\Psi D O^{\prime}$ instead of "pseudo differential operators".

The next theorem states that $p(x, D)$ is a continuous linear map of $C_{o}^{\infty}(\Omega)$ into $C^{\infty}\left(\Omega\right.$ which extends to $E^{\prime}(\Omega)$. For the proof, we need a result which depends on

Lemma 4.5. Let $p \in S^{m}(\Omega)$ and $\phi \epsilon C_{o}^{\infty}(\Omega)$. Then, for each positive integer $N$, there exists $c_{N}>0$ such that for all $\xi, \eta$ in $\mathbb{R}^{n}$,

$$
\left|\int e^{2 \pi i x \cdot \eta} p(x, \xi) \phi(x) d x\right| \leq c_{N}(1+|\xi|)^{m}(1+|\eta|)^{-N}
$$

Proof. For any $\xi$ and $\eta$ in $\mathbb{R}^{n}$,
$\left|\eta^{\alpha} \int e^{2 \pi i x . \eta} p(x, \xi) \phi(x) d x\right|$

$$
=\left|\int D_{x}^{\alpha} e^{2 \pi i x \cdot \eta} p(x, \xi) \phi(x) d x\right|
$$

Integrating by parts, we have

$$
\begin{aligned}
\left|\eta^{\alpha} \int e^{2 \pi i x . \xi} p(x, \xi) \phi(x) d x\right| & =\left|\int e^{2 \pi i x . \eta} D_{x}^{\alpha}(p(x, \xi) \phi(x)) d x\right| \\
& \leq C_{\alpha}(1+|\xi|)^{m} \text { for all } \alpha
\end{aligned}
$$

Therefore

$$
\sum_{|\alpha| \leq N}\left|\eta^{\alpha} \int e^{2 \pi i x . \eta} p(x, \xi) \phi(x) d x\right| \leq c_{N}^{\prime}(1+|\xi|)^{m} \text { for all } N
$$

Since $(1+\mid \eta)^{N} \leq c \sum_{|\alpha| \leq N}\left|\eta^{\alpha}\right|$, the required result follows.

Corollary 4.6. If $p \epsilon S^{m}(\Omega)$ and $\phi \epsilon C_{o}^{\infty}(\Omega)$, then the function

$$
g(\xi)=\int e^{2 \pi i x \cdot \xi} p(x, \xi) \phi(x) d x
$$

is rapidly decreasing as $\xi$ tends to $\infty$.
Proof. Set $\xi=\eta$ in the lemma.
Theorem 4.7. If $p \epsilon S^{m}(\Omega)$, then $p(x, D)$ is a continuous linear map of $C_{o}^{\infty}(\Omega)$ into $C^{\infty}(\Omega)$ which can be extended as a linear map from $E^{\prime}(\Omega)$ into $D^{\prime}(\Omega)$.

Proof. For $u \epsilon C_{o}^{\infty}(\Omega)$,

$$
p(x, D) u(x)=\int e^{2 \pi i x . \xi} p(x, \xi) \hat{u}(\xi) d \xi
$$

The integral converges absolutely and uniformly on compact sets, as do the integrals

$$
\int D_{x}^{\alpha}\left(e^{2 \pi i x \cdot \xi} p(x, \xi) \hat{u}(\xi) d \xi \text { for all } \alpha\right.
$$

since $p \in S^{m}$ and $\hat{u} \in \mathcal{S}$.
This proves that $p(x, D) u \epsilon C^{\infty}(\Omega)$, and continuity of $p(x, D)$ from $C_{o}^{\infty}(\Omega)$ to $C^{\infty}(\Omega)$ is an easy exercise.

To prove the rest of the theorem, we will make use of Corollary 4.6
For $u \epsilon E^{\prime}(\Omega)$, we define $p(x, D) u$ as a functional on $C_{o}^{\infty}(\Omega)$ as follows:

$$
\begin{aligned}
<p(x, D) u, \phi> & =\iint p(x, \xi) \hat{u}(\xi) e^{2 \pi i x \cdot \xi} \phi(x) d x d \xi \\
& =\int g(\xi) \hat{u}(\xi) d \xi
\end{aligned}
$$

where

$$
g(\xi)=\int e^{2 \pi i x . \xi} p(x, \xi) \phi(x) d x, \text { for } \phi \epsilon C_{o}^{\infty}(\Omega)
$$

By Corollary 4.6, $g(\xi)$ is rapidly decreasing, while $\hat{u}$ is of polynomial growth; so the last integral is absolutely convergent and the functional on $C_{o}^{\infty}$ thus defined is easily seen to be continuous. Moreover, if $u \epsilon C_{o}^{\infty}$, the double integral is absolutely convergent, and by interchanging the order of integration, we see that this definition of $p(x, D) u$ coincides with the original one. Hence we have extended $p(x, D)$ to a map from $E^{\prime}(\Omega)$ to $\mathcal{D}(\Omega)$.

Remark 4.8. It follows easily from the above argument that $p(x, D)$ is sequentially from $E^{\prime}(\Omega)$ to $\mathcal{D}(\Omega)$, i.e., if $u_{k}$ converges to $u$ in $E^{\prime}(\Omega)$, then $p(x, D) u_{k}$ converges to $p(x, D) u$ in $\mathcal{D}^{\prime}(\Omega)$. Actually, $p(x, D)$ is continuous from $E^{\prime}(\Omega)$ to $\mathcal{D}^{\prime}(\Omega)$ : this follows from the fact (which we shall prove later) that the transpose of a pseudo differential operator is again a pseudo differential operator. Thus $p(x, D)^{t}: C_{o}^{\infty}(\Omega) \rightarrow C^{\infty}(\Omega)$ is continuous, and so, by duality, $p(x, D)=\left(p(x, D)^{t}\right)^{t}: E^{\prime}(\Omega) \rightarrow \mathcal{D}^{\prime}(\Omega)$ is continuous.

## 2 Distribution Kernels and the Pseudo Local Property

Definition 4.9. Let $T$ be an operator from $C_{o}^{\infty}(\Omega)$ to $C^{\infty}(\Omega)$. If there exists a distribution $K$ on $\Omega \times \Omega$ such that

$$
<T u, v>=<k, v \otimes u>\text { for } u, v \epsilon C_{o}^{\infty}(\Omega),
$$

we say that $K$ is the distribution kernel of the operator $T$.
In this definition, $v \otimes u$ is defined by $(v \otimes u)(x, y)=v(x) u(y)$. Formally, this definition says that.

$$
T u(x)=\int K(x . y) u(y) d y
$$

$K$ is uniquely determined since linear combinations of functions of the form $v \otimes u$ are dense in $C_{o}^{\infty}(\Omega \times \Omega)$.

If $p \epsilon S^{m}(\Omega)$, it is easy to compute the distribution kernel of $p(x, D)$. In fact,

$$
<p(x, D) u, v>=\iint p(x, \xi) \hat{u}(\xi) e^{2 \pi i x \cdot \xi} v(x) d \xi d x
$$

$$
=\iint p(x, \xi) e^{2 \pi i x \cdot \xi}(v \otimes u)_{2}(x, \xi) d \xi d x
$$

where $(v \otimes u)_{2}(x, \xi)$ means the Fourier transform in the second variable.
It follows immediately from the definition of $K$ that

$$
\begin{aligned}
<K, w> & =\iint e^{2 \pi i x . \xi} p(x, \xi) \hat{w_{2}}(x, \xi) d \xi d x, \forall w \in C_{o}^{\infty}(\Omega \times \Omega) . \\
\text { or }\langle K, w\rangle & =\iiint e^{2 \pi i(x-y) \cdot \xi} p(x, \xi) w(x, y) d y d \xi d x .
\end{aligned}
$$

From this it is easy to see that $K(x, y)=p_{2}^{v}(x, x-y)$ where $p_{2}^{v}(x,$.$) is$ the inverse Fourier transform of the tempered distribution $p(x$, .). In particular, this shows that $p$ is uniquely determined by the operator $p(x, D)$. If $P \epsilon \Psi^{m}(\Omega)$, we shall sometimes denote the unique $p \epsilon S^{m}(\Omega)$ such that $p=p(x, D)$ by $\sigma_{p}$.

The following theorem gives precise results on the kernel $K$ of $p(x, D)$.

Theorem 4.10. The distribution kernel $K$ of $p(x, D)$ with $p \in S^{m}(\Omega)$ is in $C^{\infty}$ on $(\Omega \times \Omega) / \Delta$ where $\Delta\{(x, x): x \in \Omega\}$ is the diagonal. More precisely, if $|\alpha|>M+n+j$ for a positive integer $j$, then $(x-y)^{\alpha} K(x, y) \epsilon C^{j}(\Omega \times \Omega)$.
$88 \quad$ Proof. For $w \in C_{o}^{\infty}(\Omega \times \Omega)$, let us compute $\left\langle(x-y)^{\alpha} K, w\right\rangle$.

$$
\begin{aligned}
<(x-y)^{\alpha} K, w> & =<K,(x-y)^{\alpha} w> \\
& =\iint e^{2 \pi i x \cdot \xi} p(x, \xi)\left(x+D_{\xi}\right)^{\alpha} \hat{w_{2}}(x, \xi) d \xi d x \\
& =\iint \hat{w_{2}}(x, \xi)\left(x-D_{\xi}\right)^{\alpha}\left\{e^{2 \pi i x \cdot \xi} p(x, \xi)\right\} d \xi d x
\end{aligned}
$$

Using the Leibniz formula, it is easily seen that

$$
\left(x-D_{\xi}\right)^{\alpha}\left\{e^{2 \pi i x . \xi} p(x, \xi)\right\}=\left(-D_{\xi}\right)^{\alpha} p(x, \xi) \cdot e^{2 \pi i x \cdot \xi}
$$

Therefore

$$
\begin{aligned}
<(x-y)^{\alpha} K, w> & =\iint \hat{w_{2}}(x, \xi) e^{2 \pi i x \cdot \xi}\left(-D_{\xi}\right)^{\alpha} p(x, \xi) d \xi d x \\
& =\iiint w(x, y) e^{2 \pi i(x-y) \cdot \xi}\left(-D_{\xi}\right)^{\alpha} p(x, \xi) d y d \xi d x
\end{aligned}
$$

From the above expression, we infer that

$$
(x-y)^{\alpha} K(x, y)=\int e^{2 \pi i(x-y) \cdot \xi}\left(-D_{\xi}\right)^{\alpha} p(x, \xi) d \xi
$$

Since $\left|\left(-D_{\xi}\right)^{\alpha} p(x, \xi)\right| \leq c(1+|\xi|)^{m-|\alpha|}$, the integral converges absolutely and uniformly on compact sets whenever $m-|\alpha|<-n$. Also we can differentiate with respect to $x$ and $y \quad j$ times under the integral sign provided $m-|\alpha|+j<-n$ or $|\alpha|>m+n+j$. Thus, we see that $(x-y)^{\alpha} K \epsilon C^{j}(\Omega \times \Omega)$.

Corollary 4.11 (PSEUDO LOCAL PROPERTY OF $\Psi D O$ ). If $P \epsilon \Psi^{m}$ $(\Omega)$, then, for all $u \in E^{\prime}(\Omega)$, sing supp $P u$ is contained in sing supp $u$.

Proof. Let $u \epsilon E^{\prime}$ and let $V$ be an arbitrary neighbourhood of sing supp $u$. Take a $\phi \epsilon C_{o}^{\infty}(V)$ such that $\phi=1$ on sing supp $u$. Then $u=\phi u+(1-\phi) u=$ $u_{1}+u_{2} ; u_{2}$ is a $C_{o}^{\infty}$ function and supp $u_{1} \epsilon V$.

Therefore, $P u=P u_{1}+P u_{2}$ and $P u_{2}$ is a $C^{\infty}$ function. Moreover, when $x_{o} \notin V$,

$$
P u_{1}(x)=\int_{V} K(x, y) u_{1}(y) d y
$$

is also a $C^{\infty}$ function in a neighbourhood of $x_{o}$ since $u_{1}(y)=0$ for $y$ near $x_{o}$. This implies that sing supp $P u \subset V$. Since $V$ is any neighbourhood of sing supp $u$, the corollary is proved.

Corollary 4.12. If $p \in S^{-\infty}(\Omega)$, then the distribution kernel $K$ of $p(x, D)$ is in $C^{\infty}(\Omega \times \Omega)$.

Proof. Follows from the theorem if we take $|\alpha|=0$.
Corollary 4.13. If $p \in S^{-\infty}(\Omega)$, then $p(x, D)$, maps $E^{\prime}(\Omega)$ continuously into $C^{\infty}(\Omega)$.

Proof. If $u \epsilon E^{\prime}(\Omega)$, in view of Corollary 4.12, it is easily seen that $p(x, D) u$ is a smooth function defined by

$$
P(x, D) u(x)=\int K(x, y) u(y) d y=<u, K(x, .)>
$$

whence the result follows.

Definition 4.14. A smoothing operator is a linear operator $T$ which maps $E^{\prime}(\Omega)$ continuously into $C^{\infty}(\Omega)$.

If $K \epsilon C^{\infty}(\Omega \times \Omega)$, then the operator $T$ defined by

$$
(T f)(x)=<K(x, .), f\rangle=\int K(x, y) f(y) d y
$$

is a smoothing operator. Conversely, every smoothing operator $T$ can be given in the above form with $K(x, y)=\left(T \delta_{y}\right)(x)$.

As we have already remarked if $p \epsilon S^{-\infty}$, then the corresponding $\Psi D O$ is smoothing. However, not every smoothing operator is a $\Psi D O$. For example, we have

90 Proposition 4.15. Suppose $p\left(x\right.$, .) is a $C^{\infty}$ function of $x$ with values in $E^{\prime}$. Then $p(x, D)$ (defined in the same way as in the case of $p \epsilon S^{m}(\Omega)$ ) is a smoothing operator.

Proof. The distribution kernel $K$ of $p(x, D)$ is given by

$$
\begin{aligned}
K(x, y) & =\int e^{2 \pi i(x-y) \cdot \xi} p(x, \xi) d \xi \\
& =\left\langle p(x, .), e^{2 \pi i(x-y)(.)}\right\rangle
\end{aligned}
$$

and hence $K(x, y) \epsilon C^{\infty}$. Thus $K$ defines a smoothing operator.
Remark 4.16. Sometimes, it is convenient to enlarge the class of pseudo-differential operators of order $m$ by including operators of the form $P+S$ where $P \epsilon \Psi^{m}(\Omega)$ and $S$ is smoothing. However, the general philosophy is the following:

1. On the level of operators, smoothing operators are negligible.
2. On the level of symbols, what counts is the asymptotic behaviour at $\infty$, so that symbols of order $-\infty$ are negligible.

## 3 Asymptotic Expansions of Symbols

Definition 4.17. Suppose $m_{0}>m_{1}>m_{2}>\cdots m_{j} \in \mathbb{R}, m_{j}$ tends to $-\infty$, and $p_{j} \epsilon S^{m_{j}}(\Omega), p \epsilon S^{m_{o}}(\Omega)$. We say that $p \sim \sum_{j=0}^{\infty} p_{j}$ if $p-\sum_{j<k} p_{j} \epsilon S^{m_{k}}(\Omega)$ for all $k$.

Proposition 4.18. Suppose $m_{j}$ tends to $-\infty$ and $p_{j} \epsilon S^{m j}(\Omega)$. Then there exists a $p$ in $S^{m_{o}}(\Omega)$ such that $p \sim \sum^{\infty} j=0 p_{j}$. This $p$ is unique modulo $S^{-\infty}(\Omega)$.

Proof. Let $\left(\Omega_{n}\right)$ be an increasing sequence of compact subsets of $\Omega$ whose union is $\Omega$. Fix $\phi \epsilon C^{\infty}$ with $\phi=1$ for $|\xi| \geq 1>$ and $\phi=0$ for $|\xi| \leq \frac{1}{2}$.

Claim There exists a sequence $\left(t_{j}\right), t_{j} \geq 0$ and $t_{j}$ tending to $\infty$ so 91 rapidly that we have
(4.19)
$\left|D_{x}^{\beta} D_{\xi}^{\alpha}\left(\phi\left(\xi / t_{j}\right) p_{j}(x, \xi)\right)\right| \leq 2^{-j}(1+|\xi|)^{m_{j-1}-|\alpha|}$ for $x \epsilon \Omega_{i}$, and $|\alpha|+|\beta|+i \leq j$.
Granted this, we define

$$
p(x, \xi)=\sum_{j=0}^{\infty} \phi\left(\xi / t_{j}\right) p_{j}(x, \xi)
$$

Note that for each $x$ and $\xi$, the sum is finite. Using (4.19), it is straightforward to show that $p \epsilon S^{m_{o}}(\Omega)$ and $p \sim \sum_{j=0}^{\infty} p_{j}$. Moreover, suppose that $q \epsilon S^{m_{o}}(\Omega)$ and $q \sim \sum_{j=0}^{\infty} p_{j}$. Then

$$
p-q=\left(p-\sum_{j<k} p_{j}\right)-\left(q-\sum_{j<k} p_{j}\right) \epsilon S^{m_{k}}(\Omega) \text { for all } k .
$$

Hence $p-q \epsilon S^{-\infty}$.
we now briefly indicate the steps involved in proving the claim :
i) First observe that $\left|D_{\xi}^{\nu} \phi\left(\xi / t_{j}\right)\right| \leq c|\xi|^{-|\nu|}$ uniformly in $j$.
ii) Hence we have $\left|D_{x}^{\beta} D_{\xi}^{\alpha}\left(\phi\left(\xi / t_{j}\right) p_{j}(x, \xi)\right)\right| \leq c_{j}(1+|\xi|)^{m_{j}-|\alpha|}$ for $x \in \Omega_{i}$, $|\alpha|+|\beta|+i \leq j$.
iii) Finally, pick $t_{j}$ so large that " $|\xi| \geq t_{j} / 2 \geq c_{j}(1+|\xi|)^{m_{j}-m_{j-1}} \leq 2^{-j}$, Details are left to the reader.

The following theorem provides a useful criterion for the asymptotic relation $p \sim \sum_{j=0}^{\infty} p_{j}$ to hold.

Theorem 4.20. Suppose $p_{j} \epsilon S^{m_{j}}, m_{j}$ tends to $-\infty$, and $p \in C^{\infty}\left(\Omega \times \mathbb{R}^{n}\right)$ satisfies following conditions:
i) for all $\alpha$ and $\beta$ and all $\Omega^{\prime} \Subset \Omega$, there exists $c>0, \mu \in \mathbb{R}$ such that

$$
\left|D_{x}^{\beta} D_{\xi}^{\alpha} p(x, \xi)\right| \leq c(1+|\xi|)^{\mu}, x \in \Omega, \text { and }
$$

ii) there exists a sequence $\left(\mu_{k}\right), \mu_{k}$ tending to $-\infty$ so that $\mid p(x, \xi)-$ $\sum_{j<k} p_{j}(x, \xi) \mid \leq C_{\Omega},(1+|\xi|)^{\mu_{k}}$ for $x \in \Omega^{\prime}$. Then $p \in S^{m_{o}}(\Omega)$ and $p \sim$ $\sum_{j=0}^{\infty} p_{j}$.

To prove this theorem, we need the following
Lemma 4.21. Let $\Omega_{1}$ and $\Omega_{2}$ be two compact subsets of $\mathbb{R}^{n}$ such that $\Omega_{1}$ is contained in the interior of $\Omega_{2}$. Then there exists constants $c_{1}>0$ and $c_{2}>0$ such that for all $f \in C^{2}\left(\Omega_{2}\right)$,

$$
\operatorname{Sup}_{\Omega_{1}}\left|\partial_{j} f\right|^{2} \leq c_{1}\left(\operatorname{Sup}_{\Omega_{2}}|f|^{2}\right)+c_{2}\left(\operatorname{Sup}_{\Omega_{2}}|f|\right)\left(\operatorname{Sup}_{\Omega_{2}}\left|\partial_{j}^{2} f\right|\right) .
$$

Proof. It suffices to assume that $n=1$ and $f$ is real valued. With this reduction, the proof becomes an exercise in elementary calculus. This idea is roughly as follows: We wish to show that if $\|f\|_{\infty}$ and $\left\|f^{\prime \prime}\right\|_{\infty}$ are both small, then so is $\left\|f^{\prime}\right\|_{\infty}$ is small and $\left|f^{\prime}\left(x_{o}\right)\right|$ is large, then $\left|f^{\prime}\right|$ will be large in some sizable interval $[a, b]$ containing $x_{o}$. But then $|f(b)-f(a)|$ and hence $\|f\|_{\infty}$ is large. We leave it to the reader to work out the quantitative details of this argument.

Proof of the theorem. By Proposition 4.18 there exists $q$ in $S^{m_{o}}$ such that $q \sim \sum_{j=0}^{\infty} p_{j}$; so, it will suffice to show that $p-q \epsilon S^{-\infty}$. First,

$$
|p(x, \xi)-q(x, \xi)|=\left|\left(p(x, \xi)-\sum_{j<k} p_{j}(x, \xi)\right)-\left(q(x, \xi)-\sum_{j<k} p_{j}(x, \xi)\right)\right|
$$

For

$$
\Omega^{\prime} \subset \subset \Omega,\left|p(x, \xi)-\sum_{j<k} p_{j}(x, \xi)\right| \leq c_{\Omega},(1+|\xi|)^{\mu_{k}}, x \in \Omega^{\prime}
$$

Also $q \sim \sum_{j=0}^{\infty} p_{j}$. These show that, for any $N$,

$$
|p(x, \xi)|-q(x, \xi) \mid \leq C_{N \Omega^{\prime}}(1+|\xi|)^{-N}, x \epsilon \Omega^{\prime}
$$

We want to prove such an estimate for $D_{x}^{\beta} D_{\xi}^{\alpha}(p-q)$ also. To this end, we will apply Lemma 4.21 to the function $(x, \eta) \rightarrow(p-q)(x, \xi+$ $\eta$ ) considering $\xi$ as a parameter, and taking $\Omega_{1}=\Omega^{\prime} \times 0$ and $\Omega_{2}$ a small neighbourhood of $\Omega_{1}$. If $|\alpha|+|\beta|=1$, we use the estimate just established for $p-q$, together with the hypothesis (i) on the second order derivatives of $p$; the lemma implies that $D_{x}^{\beta} D_{\xi}^{\alpha}(p-q)$ is rapidly decreasing for $|\alpha|+|\beta|=1$. Combining this with the hypothesis (i) on the third order derivatives of $p$, we see that $D_{x}^{\beta} D_{\xi}^{\alpha}$ is rapidly decreasing for $|\alpha|+|\beta|=2$. Proceeding by induction on $|\alpha|+|\beta|$ we get the required result.

## 4 Properly Supported Operators

Since pseudo differential operators map $C_{o}^{\infty}$ to $C^{\infty}$ rather then $C_{o}^{\infty}$, it is generally not possible to compose two of them. The problem may be remedied by considering a more restricted class of operators, the so called "properly supported" ones.

Definition 4.22. A subset $K$ of $\Omega \times \Omega$ is said to be proper if, for any $\mathbf{9 4}$ compact set $\Omega^{\prime} \subset \Omega$, both $\pi_{1}^{-1}\left(\Omega^{\prime}\right) \cap K$ and $\pi_{2}^{-1}\left(\Omega^{\prime}\right) \cap K$ are compact.

Here $\pi_{1}$ and $\pi_{2}$ are projections of $\Omega \times \Omega$ onto the first and second factors.

For example the diagonal $\Delta=\{(x, x): x \in \Omega\}$ is proper. Most of the proper subsets we will be considering are neighbourhoods of subsets of the diagonal.


Figure 4.1: A proper set
Definition 4.23. An operator $T: C_{o}^{\infty}(\Omega) \rightarrow C^{\infty}(\Omega)$ is said to be properly supported, if its distribution kernel $K$ has proper support.
Exercise. Let $T=\sum a_{\alpha}(x) D^{\alpha}$ be a differential operator on $\Omega$. Computer the distribution kernel $K$ of $T$ and show that supp $K$ is a subset of the diagonal. Thus $T$ is properly supported.

If $T$ is properly supported then it maps $C_{o}^{\infty}$ into itself, since supp $T u \subset \pi_{1}\left(\pi_{2}^{-1}(\operatorname{supp} u) \cap \operatorname{supp} K\right)$, as is easily seen from the formula $T u(x)=\int K(x, y) u(y) d y$. More generally, for any $\Omega^{\prime} \subset \subset \Omega$, there exists $\Omega^{\prime \prime} \subset \subset \Omega$ such that the values of $T u$ on $\Omega^{\prime}$ depend only on the values of $u$ on $\Omega^{\prime \prime}$ namely, $\Omega^{\prime \prime}=\pi_{2}\left(\pi^{-1}\left(\Omega^{\prime}\right) \cap \operatorname{supp} K\right)$. From this it follows that $T$ can be extended to a map from $C^{\infty}(\Omega)$ to itself. In fact, if $u \epsilon C^{\infty}(\Omega)$ and $\Omega^{\prime}, \Omega^{\prime \prime}$ are as above, we define Tu on $\Omega^{\prime}$ by

$$
\left.T u\right|_{\Omega^{\prime}}=\left.T(\phi u)\right|_{\Omega},
$$

where $\phi \epsilon C^{\infty}(\Omega)$ and $\phi=1$ on $\Omega^{\prime \prime}$. This definition is independent of the choice of $\phi$ and gives the same result on the intersection of two $\Omega^{\prime} s$; so $T u$ is well defined on all of $\Omega$.

If $T$ is a properly supported pseudo differential operator, so that $T$ extends to a map from $E^{\prime}(\Omega)$ to $\mathcal{D}^{\prime}(\Omega)$, the same arguments show that $T$ maps $E^{\prime}(\Omega)$ into itself and extends further to a map of $\mathcal{D}^{\prime}(\Omega)$ to itself.

Suppose $T$ and $S$ are properly supported operators on $C^{\infty}(\Omega)$ with distribution kernels $K$ and $L$ which are $C^{\infty}$ off the diagonal. Then $T S$
is an operator on $C^{\infty}(\Omega)$ with distribution kernel $M$ formally given by $M(x, y)=\int K(x, z) L(z, y) d z$. In fact, if $x \neq y$, since $K(x,$.$) and L(., Y)$ are smooth except at $x$ and $y$, the product $K(x,) L.(., Y)$ is a well defined element of $E^{\prime}$, and the formula $M(x, y)=<K(x,) L.(., y), 1>$ displays $M$ as a $C^{\infty}$ function off the diagonal.

Proposition 4.24. Supp $M$ is proper. Thus, TS is properly supported.
Proof. Clearly,
Supp $M \subset\left\{(x, y): \pi_{2}\left(\pi_{1}^{-1}(x) \cap \operatorname{supp} K\right) \cap \pi_{1}\left(\pi_{2}^{-1}(y) \cap \operatorname{supp} L\right) \neq \phi\right\}$ Suppose $A \subset \Omega$ is compact and set $B=\pi_{2}\left(\pi_{1}^{-1}(A) \cap \operatorname{supp} K\right)$, then $B$ is compact and

$$
\begin{aligned}
\pi_{1}^{-1}(A) \cap \operatorname{supp} M & \subset A \times\left\{y: B \cap \pi_{1}\left(\pi_{2}^{-1}(y) \cap \operatorname{supp} L\right) \neq \phi\right\} \\
& =A \times\left\{y: \pi_{1}^{-1}(B) \cap \pi_{2}^{-1}(y) \cap \operatorname{supp} L \neq \phi\right\} \\
& =A \times \pi_{2}\left(\pi_{1}^{-1}(B) \cap \operatorname{supp} L\right)
\end{aligned}
$$

which is compact. Likewise $\pi_{2}^{-1}(A) \cap \operatorname{supp} M$ is also compact.
Exercise. Suppose $A \subset(\Omega \times \Omega)$ is proper. Show that there exists a properly supported $\phi \epsilon C^{\infty}(\Omega \times \Omega)$ such that $\phi=1$ on $A$. (Hint: Let $\left\{\phi_{j}\right\} \subset$ $C_{o}^{\infty}(\Omega \times \Omega)$ be a partition of unity on $\Omega \times \Omega$ and let $\left.\phi=\sum_{A \cap s u p p} \phi_{j}=\phi^{\phi_{j}}\right)$.

## $5 \psi d o^{\prime} s$ Defined by Multiple Symbols

We have arranged our definition of pseudo differential operators to agree with the usual convention for differential operators, according to which differentiations are performed first, followed by multiplication by the coefficients. However, in some situations (for example, computing adjoints), it is convenient to have a more flexible setup which allows multiplication operators both before and after differentiations. We therefore, introduce the following apparently more general class of operators (which, however, turns out to coincide with the class of $\psi D O$, modulo smoothing operators).

Definition 4.25. For $\Omega$ open in $\mathbb{R}^{n}$ and $m$ real, we define the class of multiple symbols of order $m$ on $\Omega$,

$$
S^{m}(\Omega \times \Omega)=\left\{a \epsilon C^{\infty}\left(\Omega \times \mathbb{R}^{n} \times \Omega\right): \forall \alpha, \beta, \forall \Omega^{\prime} \subset \subset \Omega, \exists c=c_{\alpha \beta \Omega^{\prime}}\right.
$$

such that $\left.\sup _{x, y \in \Omega^{\prime}}\left|D_{x, y}^{\beta} D_{\xi}^{\alpha} a(x, \xi, y)\right| \leq c(1+|\xi|)^{m-|\alpha|}\right\}$
When $a \epsilon S^{m}(\Omega \times \Omega)$ we define the operator $A=a(x, D, y)$ by

$$
A u(x)=\iint e^{2 \pi i(x-y) \cdot \xi} a(x, \xi, y) u(y) d y d \xi
$$

Here the integral must be interpreted as an iterated integral with integration performed first in $y$ then in $\xi$ as it is not absolutely convergent as a double integral.

We observe that if $a(x, \xi, y)=a(x, \xi)$ is independent of $y$, then $A=a(x, D)$; thus this class of operators include the $\psi D O^{\prime} s$. We also observe that different $a^{\prime} s$ may give rise to same operator. For example, if $a(x, \xi, y)=\phi(x) \psi(y)$ with $\psi, \phi \epsilon C_{o}^{\infty}(\Omega)$ and $\operatorname{supp} \psi \cap \operatorname{supp} \phi=\phi$, then $a \epsilon s^{o}(\Omega \times \Omega)$ and $a(x, D, y)=0$.

Definition 4.26. Given $a \epsilon S^{m}(\Omega \times \Omega)$ we define

$$
\sum_{a}=\left\{(x, y) \epsilon \Omega \times \Omega:(x, \xi, y) \epsilon \text { supp a for some } \xi \in \mathbb{R}^{n}\right\}
$$

Proposition 4.27. Let $a \epsilon S^{m}(\Omega \times \Omega)$ and let $K$ be the distribution kernel of $A=a(x, D, y)$. Then
i) $\operatorname{Supp} K \subset \sum_{a}$, and
ii) if the support of $K$ is proper, then there exists $a^{\prime} \epsilon S^{m}(\Omega \times \Omega)$ such that $a^{\prime}(x, \xi, y)=a(x, \xi, y)$ when $(x, y)$ is near the diagonal in $\Omega \times \Omega, \sum_{a}$, is proper and $a(x, D, y)=a^{\prime}(x, D, y)$.

Proof. The kernel $K$ is given by

$$
<K, w>=\iiint e^{2 \pi i(x-y) \cdot \xi} a(x, \xi, y) w(x, y) d y d \xi d x
$$

From this, we see that $\langle K, w\rangle .=0$ whenever supp $w \cap \sum_{a}=\phi$. Therefore, $\operatorname{supp} K \subset \sum_{a}$. This proves (i).

To prove (ii), choose a properly supported $\phi \epsilon C^{\infty}(\Omega \times \Omega)$ with $\phi=1$ on $\Delta \cup \operatorname{supp} K, \Delta$ being the diagonal. Set $a(x, \xi, y)=\phi(x, y) a(x, \xi, y)$. Then $\sum_{a}, \subset \operatorname{supp} \phi$ and hence $\sum_{a}$, is proper. Also $a^{\prime}(x, \xi, y)=a(x, \xi, y)$ when $(x, y)$ is near the diagonal. Now,

$$
\begin{aligned}
<a(x, D, y) u, v> & =<K, v \otimes u> \\
& =<\phi K, v \otimes u> \\
& =<K, \phi(v \otimes u)> \\
& =\iiint e^{2 \pi i(x-y) \cdot \xi} a(x, \xi, y) \phi(x, y) \times v(x) u(y) d y d \xi d x \\
& =\iiint e^{2 \pi i(x-y) \cdot \xi} a^{\prime}(x, \xi, y) u(y) d y d \xi d x \\
& =<a^{\prime}(x, D, y) u, v>
\end{aligned}
$$

Since $v$ is arbitrary, we have $a(x, D, y)=a^{\prime}(x, D, y)$.
Theorem 4.28. Suppose $a \epsilon S^{m}(\Omega \times \Omega)$ and $A=a(x, D, y)$ is properly supported. Let $p(x, \xi)=e^{-2 \pi i x \cdot \xi} A\left(e^{2 \pi i(.) \cdot \xi}\right)(x)$. Then $p \in S^{m}(\Omega)$ and $p(x, D)=A$. Further,

$$
\left.p(x, \xi) \sim \sum_{a} \frac{1}{\alpha!} \partial_{\xi}^{\alpha} D_{y}^{\alpha} a(x, \xi, y)\right|_{y=x}
$$

Proof. For u in $C_{o}^{\infty}(\Omega)$, we have

$$
u(x)=\int e^{2 \pi i x \cdot \xi} \hat{u}(\xi) d \xi
$$

Therefore, by linearity and continuity of $A$,

$$
\begin{aligned}
A u(x) & =\int A\left(e^{2 \pi i(\cdot) \cdot \xi}\right)(x) \hat{u}(\xi) d \xi \\
& =\int e^{2 \pi i x \cdot \xi} p(x, \xi) \hat{u}(\xi) d \xi=p(x, D) u(x)
\end{aligned}
$$

By the proposition above, we can modify a so that $\sum_{a}$ is proper without affecting $A$ or the behaviour of a along the diagonal $x=y$ (which is what enters into the asymptotic expansion of $p$ ). Set $b(x, \xi, y)=$ $a(x, \xi, x+y)$. Then as a function of $y, b$ has compact support. Indeed, for fixed $x$ and $\xi$, if $y \in\left\{y^{\prime}: b\left(x, \xi, y^{\prime}\right) \neq 0\right\}$ then $x+y \epsilon \pi_{2}\left(\pi_{1}^{-1}(x) \cap \sum_{a}\right)$ which is compact. Now,

$$
\begin{aligned}
p(x, \eta) & =e^{-2 \pi i x \cdot \eta} \iint e^{2 \pi i(x-y) \cdot \xi} b(x, \xi, y-x) e^{2 \pi i y \cdot \eta} d y d \xi \\
& =e^{-2 \pi i x \cdot \eta} \iint e^{2 \pi i z \cdot \xi} b(x, \xi z) e^{2 \pi i(z+x) \cdot \eta} d z d \xi \\
& =\int \hat{b}_{3}(x, \xi, \xi-\eta) d \xi=\int \hat{b}_{3}(x, \xi+\eta, \xi) d \xi
\end{aligned}
$$

where $\hat{b}_{3}$ is the Fourier transform of $b$ in the third variable. Since $b(x, \xi,$.$) is in C_{o}^{\infty}$, this calculation is justified and $\hat{b}_{3}(x, \eta, \xi)$ is rapidly decreasing in the variable $\xi$.

More precisely, since $a \epsilon S^{m}(\Omega \times \Omega)$, we have

$$
\begin{equation*}
\left|D_{x}^{\beta} D_{\eta}^{\alpha} \hat{b}_{3}(x, \eta, \xi)\right| \leq c_{\alpha \beta N}(1+|\eta|)^{m-|\alpha|}(1+|\xi|)^{-N} \tag{4.29}
\end{equation*}
$$

Since

$$
\left(\frac{1+|\eta|}{1+|\xi|}\right)^{s} \leq(1+|\xi-\eta|)^{|s|}, \text { taking } \eta+\xi \text { instead of } \eta
$$

we have $(1+|\xi+\eta|)^{s} \leq(1+|\xi|)^{s}(1+|\eta|)^{|s|}$. If we use this in

$$
\left|D_{x}^{\beta} D_{\eta}^{\alpha} \hat{b}_{3}(x, \xi+\eta, \xi)\right| \leq c_{\alpha \beta N}(1+|\xi+\eta|)^{m-|\alpha|}(1+|\xi|)^{-N}
$$

we get

$$
\left|D_{x}^{\beta} D_{\eta}^{\alpha} \hat{b}_{3}(x, \xi+\eta, \xi)\right| \leq c_{\alpha \beta N}^{\prime}(1+|\xi|)^{-N+m-|\alpha|}(1+|\eta|)^{|m|+|\alpha|}
$$

If we take $N$ so that $-N+m-|\alpha|<-n$ we have

$$
\begin{align*}
\left|D_{x}^{\beta} D_{\eta}^{\alpha} p(x, \eta)\right| & =\left|\int D_{x}^{\beta} D_{\eta}^{\alpha} \hat{b}_{3}(x, \xi+\eta, \xi) d \xi\right|  \tag{4.30}\\
& \leq c_{\alpha \beta}(1+|\eta|)^{|m|+|\alpha|}
\end{align*}
$$

100 On the other hand, taking the Taylor expansion of $\hat{b}_{3}(x, \xi+\eta, \xi)$ in the middle argument about the point $\eta$, 4.29) gives

$$
\begin{aligned}
\mid \hat{b}_{3}(x, \xi+\eta, \xi) & \left.-\sum_{|\alpha|<k} \partial_{\eta}^{\alpha} \hat{b}_{3}(x, \eta, \xi) \frac{\xi^{\alpha}}{\alpha!} \right\rvert\, \\
& \leq C|\xi|^{k} \sup _{\substack{|\alpha| \mid=k \\
0 \leq t \leq 1}}\left|\partial_{\eta}^{\alpha} \hat{b}_{3}(x, \eta+t \xi, \xi)\right| \\
& \leq C_{N}|\xi|^{k} \sup _{0 \leq t \leq 1}(1+|\eta+t \xi|)^{m-k}(1+|\xi|)^{-N}
\end{aligned}
$$

When $|\xi|<\frac{1}{2}|\eta|$, taking $N=k$ we get

$$
\left|\hat{b}_{3}(x, \xi+\eta, \xi)-\sum_{|\alpha|<k} \partial_{\eta}^{\alpha} \hat{b}_{3}(x, \eta, \xi) \frac{\xi^{\alpha}}{\alpha!}\right| \leq c_{k}(1+|\eta|)^{m-k}
$$

When $|\xi| \geq \frac{1}{2}|\eta|$, we have

$$
\left|\hat{b}_{3}(x, \xi+\eta, \xi)-\sum_{|\alpha|<k} \partial_{\eta}^{\alpha} \hat{b}_{3}(x, \eta, \xi) \frac{\xi^{\alpha}}{\alpha!}\right| \leq C_{N}(1+|\xi|)^{m+k-N}
$$

(Actually, the exponent of $(1+|\xi|)$ ) can be taken as $m-N$ when $m-k \geq 0$ and when $m-k<0)$. Also we have

$$
\begin{aligned}
\int \partial_{\eta}^{\alpha} \hat{b}_{3}(x, \eta, \xi) \xi^{\alpha} d \xi & =\left.D_{Y}^{\alpha} \int e^{2 \pi i Y \cdot \xi} \partial_{\eta}^{\alpha} \hat{b}_{3}(x, \eta, \xi) d \xi\right|_{Y=0} \\
& =\left.D_{Y}^{\alpha} \partial_{\eta}^{\alpha} b(x, \eta, y)\right|_{Y=0} \\
& =\left.D_{Y}^{\alpha} \partial_{\eta}^{\alpha} a(x, \eta, y)\right|_{y=x}
\end{aligned}
$$

So finally

$$
\begin{aligned}
\mid p(x, \eta) & \left.-\left.\sum_{|\alpha|<k} D_{y}^{\alpha} \partial_{\eta}^{\alpha} a(x, \eta, y) \frac{1}{\alpha!}\right|_{y=x} \right\rvert\, \\
& =\left|\int \hat{b}_{3}(x, \xi+\eta, \xi) d \xi-\sum_{|\alpha|<k} \frac{1}{\alpha!} \int \partial_{\eta}^{\alpha} \hat{b}_{3}(x, \eta, \xi) \xi^{\alpha} d \xi\right|
\end{aligned}
$$

$$
\leq C_{k} \int_{|\xi|<\frac{|\eta|}{2}}(1+|\eta|)^{m-k} d \xi+C_{N} \int_{|\xi| \geq \frac{|\eta|}{2}}(1+|\xi|)^{m+k-N} d \xi
$$

Now

$$
\int_{|\xi|<\left\lvert\, \frac{|\eta|}{2}\right.} C_{k}(1+|\eta|)^{m-k} d \xi=C_{k}^{\prime}(1+|\eta|)^{m-k}|\eta|^{n} \leq C_{k}^{\prime}(1+|\eta|)^{m-k+n}
$$

In the second integral, we take $N>n+m+k$ so that

$$
\int_{|\xi| \geq \frac{|\eta|}{2}} C_{N}(1+|\xi|)^{m-N} d \xi \leq C_{N}^{\prime}(1+|\eta|)^{m-N+n}
$$

by the usual integration in polar coordinates. Therefore, since $N \geq k$ we obtain

$$
\left|p(x, \eta)-\sum_{|\alpha|<k} \frac{1}{\alpha!} D_{y}^{\alpha} \sigma_{\eta}^{\alpha} a(x, \eta, y)\right|_{x=y}\left|\leq C_{k}(1+|\eta|)\right|^{\mu_{k}}
$$

with $\mu_{k}=m-k+n$.
Combining this with 4.30, we see that all the conditions of Theorem 4.20 are satisfied, so we are done.
Corollary 4.31. If $P \epsilon \psi^{m}$, there exists a $Q \epsilon \psi^{m}$ such that $Q$ is properly supported and $P-Q$ is smoothing.

Proof. If $P=p(x, D)$ choose $\phi \epsilon C^{\infty}(\Omega \times \Omega)$ which is properly supported and $\phi=1$ near the diagonal. Set $a(x, \xi, y)=\phi(x, y) p(x, \xi)$. Then $a \epsilon S^{m}(\Omega \times \Omega)$ and by construction $\sum_{a}$ is proper. So $a(x, D, y)=Q \epsilon \psi^{m}$. If $K_{p}$ and $K_{Q}$ denote the distribution kernels of $P$ and $Q$ then it is easily seen that $K_{Q}=\phi K_{p}$. This shows that $K_{Q}-K_{p}$ vanishes near the diagonal: so, by Theorem 4.10 $K_{Q}-K_{p}$ is $C^{\infty}$. Therefore, $Q-p$ is smoothing.

Remark 4.32. The idea of multiple symbols can clearly be generalised. For example, if $a \epsilon C^{\infty}\left(\Omega \times \mathbb{R}^{n} \times \Omega \times \mathbb{R}^{n}\right)$ satisfies estimates of the form $102\left|D_{x, y}^{v} D_{\xi}^{\alpha} D_{\eta}^{\beta} a(x, \xi, y, \eta)\right| \leq C(1+\eta \mid)^{m_{2}-|\beta|}(1+|\xi|)^{m_{1}-|\alpha|}$, one can define an operator $A=a(x, D, y, D)$ by

$$
A u(x)=\iiint e^{2 \pi i(x . \xi-y \cdot \xi+y \cdot \eta)} a(x, \xi, y, \eta) \hat{u}(\eta) d \eta d y d \xi
$$

which agree with our previous definitions, if a is independent of the last one or two variables. As one would expect, under reasonable restriction on supp $a$, it turns out that $A \epsilon \psi^{m_{1}+m_{2}}(\Omega)$. The reader may wish to amuse himself by working this out and computing the symbol of $A$.

## 6 Products and Adjoint of $\psi D O^{\prime} S$

We are now in a position to compute products and adjoints of $\psi D O^{\prime} S$. First we clarify our terminology. Let $T S$ be linear operators from $C_{o}^{\infty}(\Omega)$ into $C^{\infty}(\Omega)$. We say that $S=T^{t}$ if $\langle T u, v>=<u, S u\rangle$ for $u, v \in C_{o}^{\infty}(\Omega)$ and we say that $S=T^{*}$ if $\langle T u, \bar{v}\rangle=\langle u, S v\rangle$ for $u, v \in C_{o}^{\infty}(\Omega)$.

Remark 4.33. If $T$ is properly supported, then $T^{t}$ and $T^{*}$ are also properly supported. Indeed, if $K$ is the distribution kernel of $T$, then the distribution kernel of $T^{t}$ is $K^{t}$ and that of $T^{*}$ is $K^{*}$ where $K^{t}(x, y)=K(y, x)$ and $K^{*}(x, y)=K(\overline{y, x})$.

We recall that if $\operatorname{P\epsilon } \psi^{m}(\Omega)$, we denote the corresponding symbol in $S^{m}(\Omega)$ by $\sigma_{p}$.

Theorem 4.34. If $P \in \psi^{m}(\Omega)$ is properly supported, then $P^{t}, P^{*} \in \psi^{m}(\Omega)$ and $\sigma_{p} t(x, \xi) \sim \sum_{\alpha} \frac{(-1)^{|\alpha|}}{\alpha!} \partial_{\xi}^{\alpha} D_{x}^{\alpha} \sigma_{p}(x,-\xi)$,

$$
\sigma_{p^{*}}(x, \xi) \sim \sum_{a} \frac{1}{\alpha!} \partial_{\xi}^{\alpha} D_{x}^{\alpha} \bar{\sigma}_{p}(x, \xi) .
$$

Proof. For $u, v \in C_{o}^{\infty}(\Omega)$, we have

$$
\begin{aligned}
<P u, v> & =\iint e^{2 \pi i x, \xi} p(x, \xi) \hat{u}(\xi) v(x) d \xi d x \\
& =\int\left(\int e^{2 \pi i x, \xi} p(x, \xi) v(x) d x\right) \hat{u}(\xi) d \xi
\end{aligned}
$$

$$
=\int g(\xi) \hat{u}(\xi) d \xi
$$

where

$$
g(\xi)=\int e^{2 \pi i x, \xi} p(x, \xi) v(x) d x
$$

By Corollary 4.6 $g$ us a rapidly decreasing function. Thus

$$
<P u, v>=\int g \hat{u}=\int u \hat{g}=<u, P_{v}^{t}>
$$

so that

$$
\begin{align*}
P^{t} v(y)=\hat{g}(y) & =\iint e^{2 \pi i(x-y) \cdot \xi} p(x, \xi) v(x) d x d \xi  \tag{4.35}\\
& =\iint e^{2 \pi i(y-x) \cdot \xi} p(x,-\xi) v(x) d x d \xi \\
& =a(x, D, y) v(y)
\end{align*}
$$

where $a(x, \xi y)=p(y .-\xi)$. Therefore, by Theorem4.28,

$$
P^{t} \epsilon \psi^{m}(\Omega) \text { and } \sigma_{P^{t}}(x, \xi) \sim \sum_{a} \frac{(-1)^{|\alpha|}}{\alpha!} \partial_{\xi}^{\alpha} D_{x}^{\alpha} \sigma_{P}(x,-\xi)
$$

The assertions about $P^{*}$ follows along similar lines
Theorem 4.36. If $P \epsilon \psi^{m_{1}}(\Omega)$ and $Q \epsilon \psi^{m_{2}}(\Omega)$ are properly supported, then $Q P \epsilon \psi^{m_{1}+m_{2}}(\Omega)$ and

$$
\sigma_{Q P}(x, \xi) \sim \sum_{\alpha} \frac{1}{\alpha!} \partial_{\xi}^{\alpha} \sigma_{Q}(x, \xi) D_{x}^{\alpha} \sigma_{p}(x, \xi)
$$

104 Proof. Since $P=\left(P^{t}\right)^{t}$, we have

$$
(P u)(x)=\iint e^{2 \pi i(x-y) \cdot \xi} \sigma_{P}^{t}(y,-\xi) u(y) d y d \xi
$$

by 4.35.

In other words

$$
(P u)^{\wedge}(\xi)=\int e^{-2 \pi i y \cdot \xi} \sigma_{P^{t}}(y,-x i) u(y) d y
$$

Therefore,

$$
\begin{aligned}
(Q P) u(x) & =\iint e^{2 \pi i(x-y) \cdot \xi} \sigma_{Q}(x, \xi) \sigma_{P^{t}}(y,-x i) u(y) d y d \xi \\
& =a(x, D, y) u(x)
\end{aligned}
$$

where $a(x, \xi, y)=\sigma_{Q}(x, \xi) \sigma_{P^{t}}(y,-\xi)$. Clearly $a \epsilon S^{m_{1}+m_{2}}(\Omega \times \Omega)$ which implies $Q P \epsilon \psi^{m_{1}+m_{2}}(\Omega)$. Moreover

$$
\begin{aligned}
\sigma_{Q P}(x, \xi) & \left.\sim \sum_{\alpha} \frac{1}{\alpha!} \partial_{\xi}^{\alpha} D_{y}^{\alpha}\left(\sigma_{Q}(x, \xi) \sigma_{p^{t}}(y,-\xi)\right)\right|_{y=x} \\
& \sim \sum_{\alpha} \frac{1}{\alpha!} \sum_{\beta+v=\alpha} \frac{\alpha}{\beta!v!} \partial_{\xi}^{\beta} \sigma_{Q}(x, \xi) \partial_{\xi}^{v} D_{Y}^{\alpha} \sigma_{P^{t}}(x,-\xi) \\
& \sim \sum_{\alpha, v} \sum_{\delta} \frac{(-1)^{|\delta|}}{v!\beta!\delta!} \partial_{\xi}^{\beta} \sigma_{Q}(x, \xi) D^{\beta+v+\delta} \partial_{x}^{v+\delta} \partial_{\xi}^{v+\delta} \sigma_{P}(x,-\xi) \\
& \sim \sum_{\beta, \lambda}\left(\sum_{v+\delta=\lambda} \frac{(-1)^{|\delta|}}{v!\beta!\delta!}\right) \frac{1}{\beta!} \partial_{\xi}^{\beta} \sigma_{Q}(x, \xi) D_{x}^{\beta+\lambda} \partial_{\xi}^{\lambda} \sigma_{p}(x, \xi)
\end{aligned}
$$

But $\sum_{v+\delta=\lambda} \frac{(-1)^{|\delta|}}{v!\delta!}=\left(x_{0}-x_{0}\right)^{\lambda}$ with $x_{0}=(1,1, \ldots, 1)$

$$
=\left\{\begin{array}{l}
1 \text { if } \lambda=0 \\
0 \text { if } \lambda \neq 0
\end{array}\right.
$$

Therefore $\sigma_{Q P}(x, \xi) \sim \sum_{\beta} \frac{1}{\beta!} \partial_{\xi}^{\beta} \sigma_{Q}(x, \xi) D_{x}^{\beta} \sigma_{p}(x, \xi)$.

## Corollary 4.37 (PRODUCT RULES FOR $\psi D O^{\prime} S$ ).

If $Q=q(x, D) \epsilon \psi^{m}(\lambda)$ and $f \epsilon C^{\infty}(\Omega)$, then for any positive integer $N, 105$ $q(x, D)(f u)=\sum_{|\alpha|<N} \frac{1}{\alpha!} D^{\alpha} f\left[\left(\partial_{\xi}^{\alpha} q\right)(x, D) u\right]+T_{N} u$ where $T_{N} \epsilon \psi^{m-N}(\Omega)$.

Corollary 4.38. The correspondence $p \rightarrow p(x, D)$ and $P \rightarrow \sigma_{p}$ are * homomorphisms modulo lower order term i.e., if $p \in S^{m_{1}}(\Omega)$ and $q \epsilon S^{m_{2}}$ $(\Omega)$, then $(p q)(x, D)-p(x, D) q(x, D) \epsilon \psi^{m_{1}+m_{2}-1}(\Omega)$ and $p(x, D)^{*}-\bar{p}(x, D)$ $\epsilon \psi^{m_{1}-1}$ and also, if $P \epsilon \psi^{m_{1}}(\Omega)$ and $Q \epsilon \psi^{m_{2}}(\Omega)$, then $\sigma_{P Q}-\sigma_{P} \sigma_{Q}$ $\epsilon S^{m_{1}+m_{2}-1}(\Omega), \sigma_{P *}-\bar{\sigma}+p \epsilon S^{m_{1}-1}(\Omega)$.

Corollary 4.39. If $P \epsilon \psi^{m_{1}}(\Omega)$ and $Q \epsilon \psi^{m_{2}}(\Omega)$, then $[P, Q]=P Q-Q P$ $\epsilon \psi^{m_{1} m_{2}-1}(\Omega)$ and $\sigma_{[P, Q]}-\frac{1}{2 \pi i}\left\{\sigma_{P}, \sigma_{Q}\right\} \epsilon S^{m_{1} m_{2}-1}(\Omega)$ where $\{f, q\}$ stands for the Poisson bracket defined by

$$
\{f, q\}=\sum\left(\frac{\partial f}{\partial \xi_{j}} \frac{\partial g}{\partial x_{j}}-\frac{\partial f}{\partial x_{j}} \frac{\partial g}{\partial \xi_{j}}\right)
$$

Following the philosophy that smoothing operators are negligible it is a trivial matter to extend these results to non-properly supported operators.

Suppose $P \epsilon \psi^{m}(\Omega)$ is not properly supported. By corollary 4.31 we can write $P=P_{1}+s$ where $P_{1}$ is properly supported and $S$ is smoothing. Then $P^{t}=P_{1}^{t}+S$ where $P_{1}^{t}$ is given by Theorem 4.34 and $S^{t}$ is again smoothing (c.f. Remark 4.33). Likewise for $P^{*}$. If $Q$ is a properly supported $\psi D O$, the products $P Q$ and $Q P$ are well defined as operators from $C_{o}^{\infty}(\Omega)$ to $C^{\infty}(\Omega)$. Again, we have $P Q=P_{1} Q+S Q$ and $Q P=$ $Q P_{1}+Q s ; P_{1} Q$ and $Q P_{1}$ are described by Theorem 4.36, while $S Q$ and $Q S$ are smoothing.

## 7 A Continuity Theorem for $\psi$ Do on Sobolev Spaces

106 We now state and prove a continuity theorem for pseudo defferential operators acting on Sobolev spaces. We are indebted to Dr P.N. Srikanth for simplifying our original argument.

Theorem 4.40. Suppose $P=p(x, D) \epsilon \psi^{m}(\Omega)$. Then
i) $P: H_{s} \rightarrow H_{s-m}^{\mathrm{loc}}(\Omega)$ continuous for all $s \in \mathbb{R}$.
ii) If $P$ is properly supported, $P: H_{s}^{\mathrm{loc}}(\Omega) \rightarrow H_{s-m}^{\mathrm{loc}}(\Omega)$ continuously for all $s \in \mathbb{R}$

Proof. To established (i), we must show that the map $u \rightarrow \phi P u$ is bounded from $H_{s} \rightarrow H_{s-m}$ for any $\phi \epsilon C_{o}^{\infty}(\Omega)$. Replacing $p(x, \xi)$ by $\phi(x) p(x, \xi)$ we must show:
(4.41) if $p \epsilon s^{m}(\mathbb{R})$ and $p(x, \xi)=0$ for $x$ outside a compact set, then $p=p(x, D)$ is bounded from $H_{s}$ and $H_{s-m}$ for every $s$ in $\mathbb{R}$.

Suppose then that $u \epsilon H_{s}$. Then $p u \epsilon E^{\prime}$, and from the definition of $P$ on distributions, we see that

$$
\begin{aligned}
(P u)(\eta) & =<P u, e^{-2 \pi i \eta \cdot(\cdot)}> \\
& =\iint e^{2 \pi i(\xi-\eta) \cdot x} p(x, \xi) \hat{u}(\xi) d x d \xi \\
& =\int \hat{p}_{1}(\eta-\xi, \xi) \hat{u}(\xi) d \xi
\end{aligned}
$$

and hence, if $V \epsilon S$,

$$
<P u, \bar{v}>=\int(P u) \hat{\hat{v}}=\iint K(\eta, \xi) f(\xi) \bar{g}(\eta) d \xi d \eta
$$

where

$$
\begin{aligned}
& f(\xi)=\left(1+|\xi|^{2}\right)^{s / 2} \hat{u}(\xi) \\
& g(\eta)=\left(1+|\eta|^{2}\right)^{(m-s) / 2} \hat{v}(\eta)
\end{aligned}
$$

and $K(\eta, \xi)=\hat{p}_{1}(\eta-\xi, \xi)\left(1+|\xi|^{2}\right)^{-s / 2}\left(1+|\eta|^{2}\right)^{(s-m) / 2}$.
We wish to estimate $K(\eta, \xi)$. For any multi-index $\alpha$, since $p(., \xi) \epsilon C_{o}^{\infty}$, we have

$$
\begin{aligned}
\left|\zeta^{\alpha} \hat{p}_{1}(\zeta, \xi)\right| & =\mid \int\left(D_{x}^{\alpha} e^{-2 \pi i x . \zeta)} p(x, \xi) d x \mid\right. \\
& =\mid \int e^{-2 \pi i x . \zeta)}\left(D_{x}^{\alpha} p(x, \xi) d x \mid\right. \\
& \leq C_{\alpha}(1+|\xi|)^{m}
\end{aligned}
$$

so that for any positive integer $N$,

$$
\left|\hat{p}_{1}(\zeta, \xi)\right| \leq C_{N}(1+|\xi|)^{m}(1+|\zeta|)^{-N}
$$

and hence

$$
\begin{aligned}
|K(\eta, \xi)| & \leq C_{N}(1+|\xi|)^{m-s}(1+|\eta|)^{s-m}(1+|\xi-\eta|)^{N} \\
& \leq C_{N}(1+|\xi|)^{|m-s|-N}
\end{aligned}
$$

If we take $N>n+|m-s|$, we see that

$$
\int|K(\eta, \xi)| d \eta \leq C, \int|K(\eta, \xi)| d \xi \leq C .
$$

Therefore by Theorem 1.0 and the Schwarz inequality,

$$
|<P u, \bar{v}>| \leq C\|f\|_{2}\|g\|_{2}=C\|u\|_{(s)}\|v\|_{(m-s)} .
$$

From this, it follows that $\|P u\|_{(s-m)} \leq C\|u\|_{(s)}$ for $u \epsilon H_{s}$ which establishes (4.41) and hence (i).

Now suppose $P$ is properly supported and $u \epsilon H_{s}^{\text {loc }}(\Omega)$. If $\phi \epsilon C_{o}^{\infty}(\Omega)$ there exists $\Omega^{\prime} \Subset \Omega$ such that the values of Pu on $\operatorname{supp} \phi$ depend only the values of $u$ on $\Omega$. Thus if we pick $\phi^{\prime} \epsilon C_{o}^{\infty}(\Omega)$ with $\phi^{\prime}=1$ on $\Omega^{\prime}$, we have $\phi P u=\phi P\left(\phi^{\prime} u\right)$. But $\phi^{\prime} u \epsilon H_{s}$ and (4.41) $\phi P$ is bounded from $H_{s}$ to $H_{s-m}$. This establishes (ii) and completes the proof.

## 8 Elliptic Pseudo Differential Operators

Definition 4.42. $P \epsilon \psi^{m}(\Omega)$ is said to be elliptic of order $m$ if $\left|\sigma_{P}(x, \xi)\right| \geq$ $C_{\Omega},|\xi|^{m}$ for large $|\xi|$, for all $x \in \Omega^{\prime}, \Omega^{\prime} \Subset \Omega$.

Definition 4.43. If $P \epsilon \psi^{m}(\Omega), a$ left (resp. right ) parametrix of $P$ is a $\psi D O Q$ such that $Q P-I($ resp. $P Q-I)$ is smoothing.

Theorem 4.44. If $P \epsilon \psi^{m}$ is elliptic of order $m$ and properly supported, then there exists a properly supported $Q \epsilon \psi^{-m}$ which is a two-sided parametrix for $P$.

Proof. Let $P=p(x, D)$. We will obtain $Q=q(x, D)$ with $q \sim \sum_{j=0}^{\infty} q_{j}$ where the $q_{j}^{\prime}$ s are defined recursively. Let $\zeta(x, \xi)$ be a $C^{\infty}$ function such that $\zeta(x, \xi)=1$ for large $\xi$ and $\zeta(x, \xi)=0$ in a neighbourhood of the
zeros of $p$. Define $q_{o}(x, \xi)=\frac{\zeta(x, \xi)}{p(x, \xi)}$. Then $q_{o} \epsilon s^{-m}$. Let $Q_{o}=q_{o}(x, D)$.
We have

$$
\sigma_{Q_{o}^{p}}=q_{0}^{p} \quad \bmod S^{-1}=1 \quad \bmod S^{-1}=1+r_{1} \text { with } r_{1} \epsilon S^{-1}
$$

Let $q_{1}=\frac{-r_{1} \zeta}{p}, Q_{1}=q_{1}(x, D)$. Then

$$
\begin{aligned}
\sigma_{\left(Q_{0}+Q_{1}\right) P} & =\sigma_{Q_{O}^{p}}+\sigma_{Q_{1}^{p}}=1+r_{1}-r_{1} \zeta \bmod S^{-2} \\
& =1+r_{2} \text { with } r_{2} \epsilon S^{-2}
\end{aligned}
$$

Let $q_{1}=\frac{{ }^{-r} 2^{\zeta}}{p}$ etc. Having determined $q_{1}, q_{2}, \ldots q_{j}$ so that
set

$$
\begin{gathered}
\sigma_{\left(Q_{0}+Q_{1}+\cdots Q_{j}\right) P}=1+r_{j+1}, r_{j+1} \epsilon S^{-(j+1)} \\
q_{j+1}=\frac{-r_{j+1} \zeta}{p} .
\end{gathered}
$$

Let $Q$ be a properly supported operator with symbol $q \sim \sum_{j=0}^{\infty} q_{j}$. Then $\sigma_{Q P}=1 \bmod S^{-\infty}$, i.e., $Q P-I$ is smoothing. Thus $Q$ is a left parametrix. In the same way, we can construct a right parametrix $Q^{\prime}$. Then if $S=Q P-I$ and $S^{\prime}=P Q^{\prime}-I$, we have

$$
Q P Q^{\prime}=S Q^{\prime}+Q^{\prime}=Q S^{\prime}+Q
$$

So $Q-Q^{\prime}=S Q^{\prime}-Q S^{\prime}$ is smoothing. Hence $Q$ is a two sided parametrix.

Exercise. What happens if $P$ is not properly supported? (cf. the remarks at the end of \$6).

The left parametrices are used to prove regularity theorems and the right parametrices are used to prove existence theorems.

Indeed, we have:

Corollary 4.45 (ELLIPTIC REGUIARITY THEOREM). If $P$ is elliptic of order $m, u \epsilon \mathcal{D}^{\prime}(\Omega)$, $P u \epsilon H_{s}^{\text {loc }}(\Omega)$ implies $u \epsilon H_{s+m}^{\text {loc }}(\Omega)$. In particular, $P$ is hypoelliptic.

Proof. Let $Q$ be a left parametrix and set $S=Q P-I$. Then $u=Q P u-$ $S u$. Since $P u \epsilon H_{s}^{\text {loc }}(\Omega)$ and $Q$ is properly supported, $Q \epsilon \psi^{-m}$ we have $Q P u \epsilon H_{S+m}^{\mathrm{loc}}(\Omega)$ by Theorem 4.40. Also $S u \epsilon C^{\infty}$ since $S$ is smoothing. Hence $u \epsilon H_{s+m}^{\mathrm{loc}}(\Omega)$.

Theorem 4.46. Every elliptic differential operator is locally solvable. In fact, if $P$ is elliptic on $\Omega, f \in \mathcal{D}^{\prime}(\Omega)$ and $x_{o} \epsilon \Omega$, there exists $u \in \mathcal{D}^{\prime}(\Omega)$ such that $P u=f$ in a neighbourhood of $x_{o}$. (Of course, if $f \epsilon C^{\infty}$, then $u \epsilon C^{\infty}$ near $x_{o}$, by the previous corollary).

Proof. BY cutting $f$ off away from $x_{o}$, we can assume that $f \epsilon E^{\prime}$ and hence $f \epsilon H_{S}$ for some $s$. Let $Q$ be a properly supported parametrix for $P$ and let $S=P Q-I$. Then $S$ is also properly supported. If $\Omega_{1} \subset \subset \Omega$ is a neighbourhood of $x_{o}$ then there exists $\Omega_{2} \subset \subset \Omega$ such that the values of Su on $\Omega_{1}$ depend only on the values of $u$ on $\Omega_{2}$. Pick $\phi_{1}, \phi_{2} \epsilon C_{o}^{\infty}(\Omega)$ such that $\phi_{j} \equiv 1$ on $\Omega_{j}, j=1,2$ and set $T u=\phi_{1} S\left(\phi_{2} u\right)$.

Now observe the following :
i) $T u=s u$ on $\Omega_{1}$.
ii) $T: H^{s} \rightarrow C_{o}^{\infty}\left(\operatorname{suup} \phi_{1}\right)$ continuously.

From (ii) and the Arzela -Ascoli theorem, it follows that if $\left(u_{k}\right)$ is a bounded sequence in $H_{s^{\prime}}\left(T u_{k}\right)$ has a convergent subsequence in $C_{o}^{\infty}\left(\operatorname{suup} \phi_{1}\right)$ and hence in $H_{s}$. Therefore, $T$ is compact on $H_{s}$. So the equation $(T+I) u=f$ can be solved if $f$ is orthogonal (with respect to the pairing of $H_{s}$ and $\left.H_{-s}\right)$ to the space $N=\left\{g:\left(T^{*}+I\right) g=0\right\}$. This space $N$ is a finite dimensional space of $C^{\infty}$ functions so we can always make this happen by modifying $f$ outside a small neighbourhood of $x_{o}$.

Indeed, pick a basis $g_{1}, g_{2}, \ldots, g_{v}$ for $N$ and pick a neighbourhood $U$ of $x_{o}$ so small that $g_{1}, \ldots, g_{v}$ are linearly independent as functionals on $C_{o}^{\infty}(\Omega \backslash U)$. (Such a $U$ exists; otherwise, by a limiting argument using the local compactness of $N$, we could find a nontrivial linear combination of
$g_{1}, \ldots, g_{v}$ supported at $\left\{x_{o}\right\}$, which is absurd). Then we can make $f \perp N$ by adding to $f$ a function in $C_{o}^{\infty}(\Omega \backslash U)$. But then

$$
P Q u=(I+S) u=(I+T) u \text { on } \Omega_{1} \text { and }(I+T) u=f
$$

in a neighbourhood of $x_{o}$. So $Q u$ solves the problem.

## 9 Wavefront Sets

We now introduce the notion of wavefront sets, which provides a precise way of describing the singularities of distributions: it specifies not only the points at which a distribution is not smooth but the directions in which it is not smooth.

All pseudo differential operators encountered in this section will be presumed to be properly supported, and by " $\psi D O$ " we shall always mean "properly supported $\psi D O$ ".

Definition 4.47. (i) Let $\Omega$ be an open set in $\mathbb{R}^{n}$. Then we define $T^{o} \Omega=\Omega \times\left(\mathbb{R}^{n} \backslash\{0\}\right)$. (In coordinate - invariant terms, $T^{o} \Omega$ is the cotangent bundle of $\Omega$ with the zero section removed).
(ii) A set $S \subset T^{o} \Omega$ is called conic if $(x, \xi) \epsilon S \Rightarrow(x, r \xi) \epsilon S, \forall r>0$
(iii) Suppose $P=p(x, D) \epsilon \psi^{m}(\Omega)$. Then $\left(x_{o}, \xi_{o}\right) \epsilon T^{o} \Omega$ is said to be noncharacteristic for $P$ if $|p(x, \xi)| \geq c|\xi|^{m}$ for $|\xi|$ large and $(x, \xi)$ in some conic neighbourhood of $\left(x_{o}, \xi_{o}\right)$.
(iv) The characteristic variety of $P$, denoted by char $P$, is defined by char $P=\left\{(x, \xi) \epsilon T^{o} \Omega:(x, \xi)\right.$ is characteristic for $\left.p\right\}$.
(v) Let $u \in \mathcal{D}^{\prime}(\Omega)$. The wavefront set of $u$, denoted by $W F(u)$ is defined by

$$
W F(U)=\cap\left\{\operatorname{char} P: \operatorname{P\epsilon } \psi(\Omega), P u \epsilon C^{\infty}(\Omega)\right\}
$$

The restriction $P \epsilon \psi^{0}(\Omega)$ in the definition of $W F(u)$ is merely a convenient normalisation. We could allow $\psi D O$ of arbitrary order without changing anything, for if $P \epsilon \psi^{m}(\Omega)$ and $Q \epsilon \psi^{-m}(\Omega)$ is elliptic, then $Q P \epsilon \psi^{0}(\Omega)$, char $(Q P)=$ char $(P)\left(\right.$ by Theorem 4.36), and $Q P u \epsilon C^{\infty}$ if and only if $P u \epsilon C^{\infty}$ (by Corollaries 4.11 and 4.45).

Exercise. When $p=\sum_{|\alpha|=\leq m} a_{\alpha}(x) \xi^{\alpha}$ show that the characteristic variety of the differential operator $p=P(x, D)$ is

$$
\operatorname{char} P=\left\{(x, \xi): \sum_{|\alpha|=m} a_{\alpha}(x) \xi^{\alpha}=0\right\} .
$$

The motivation for $W F(u)$ is as follows. If $u \epsilon E^{\prime}$, to say that $u$ is not smooth in the direction $\xi_{o}$ should mean that $\hat{u}$ is not rapidly decreasing on the through $\xi_{o}$. On the other hand, if $P u \epsilon C_{o}^{\infty}$ then $(P u)$ should be rapidly decreasing everywhere.

If these conditions both hold, then $\xi_{o}$ must be characteristic for $P$. Localising these ideas, we arrive at $W F(u)$.

Thus $\left(x_{o} \xi_{o}\right) \epsilon W F(u)$ means roughly that $u$ fails to be $C^{\infty}$ at $x_{o}$ in the direction $\xi_{o}$. For another interpretation of this statement, see Theorem 4.56 below. For the present, we show that $W F(u)$ is related to sing $\operatorname{supp} u$ as it should be. We denote the projection $T^{o} \Omega$ onto $\Omega$ by $\pi$.

Theorem 4.48. For $u \epsilon \mathcal{D}^{\prime}(\Omega)$, sing $\operatorname{supp} u=\pi(W F(u))$.
Proof. Suppose $x_{0} \notin \sin g \operatorname{supp} u$. Then there exists $\phi \epsilon C_{o}^{\infty}$ with $\phi=1$ near $x_{o}$ and $\phi u \epsilon C_{o}^{\infty}$. But multiplication by $\phi$ is a $\psi D O$ of order 0 . Call it $P_{\phi}$. Then char $P_{\phi}=\pi^{-1}\left(\phi^{-1}(0)\right)$. Since $\phi=1$ near $x_{o}, \pi^{-1}\left(\phi^{-1}(0)\right)$ is disjoint from $\pi^{-1}\left(x_{o}\right)$. Therefore, $\pi^{-1}\left(x_{o}\right) \cap W F(u)=\phi$, i.e., $x_{o} \notin$ $\pi(W F(u))$.

Conversely, suppose $x_{o} \notin \pi(W F(u))$. Then for each $\xi$ with $|\xi|=1$, there exists $P \epsilon \psi^{0}(\Omega)$ with $P u \epsilon C^{\infty}(\Omega)$ and $\left(x_{o}, \xi\right) \notin$ char $P$.

Each char $P$ is a closed conic set; so, by the compactness of the unit sphere, there exists a finite number of $\psi D O^{\prime}$ 's say, $P_{1}, \ldots P_{N} \epsilon \psi^{0}(\Omega)$ with $P_{j} u \epsilon C^{\infty}(\Omega)$ and

$$
\left(\bigcap_{j=1}^{N} \operatorname{char} P_{j}\right) \cap \pi^{-1}\left(x_{o}\right)=\phi . \text { Set } P=\sum_{j=1}^{N} P_{j}^{*} P_{j} .
$$

Then $P$ is elliptic near $x_{o}$ and $P u \epsilon C^{\infty}(\Omega)$. Therefore by Corollary $4.45 u$ is $C^{\infty}$ near $x_{o}$, i.e., $x_{o} \notin \sin g \operatorname{supp} u$.

Definition 4.49. Let $P=p(x, D) \epsilon \psi^{m}(\Omega)$ and $U$ be an open conic set in $T^{o} \Omega$. We say that $P$ has order $-\infty$ on $U$, if, for all closed conic sets $K \subset U$ with $\pi(K)$ compact for every positive integer $N$ and multi-indices $\alpha$ and $\beta$, there exist constants $C_{\alpha \beta K N}$ such that

$$
\left|D_{x}^{\beta} D_{\xi}^{\alpha} p(x, \xi)\right| \leq C_{\alpha \beta K N}(1+|\xi|)^{-N},(x, \xi) \epsilon K
$$

The essential support of $P$ is defined to be the smallest closed conic set outside of which $P$ has order $-\infty$.

Exercise. If $P$ is a differential operator whose coefficients do not all vanish at any point of $\Omega$, show that the essential support of $P$ is $T \Omega$.

Proposition 4.50. Ess. supp $P Q \subset$ Ess. supp $P \cap$ Ess. supp $Q$.
Proof. This follows immediately from the expansion

$$
\sigma_{P Q} \sim \sum \frac{1}{\alpha!} \partial_{\xi}^{\alpha} \sigma_{P} D_{x}^{\alpha} \sigma_{C}
$$

Lemma 4.51. Let $\left(x_{o}, \xi_{o}\right) \in T^{o} \Omega$ and $U$ be any conic neighbourhood of $\left(x_{o}, \xi_{o}\right)$. Then there exists a $P \epsilon \psi^{0}(\Omega)$ such that $\left(x_{o}, \xi_{o}\right) \notin$ char $P$ but Ess. supp $P \subset U$.

Proof. Choose $p_{o}(\xi) \epsilon C^{\infty}$ with the following properties:
i) $p_{o}(\xi)=1$ for $\xi$ near $\xi_{o}, p_{o}(\xi)=0$ outside $\left\{\xi:\left(x_{o} \xi\right) \epsilon U\right\}$ and
ii) $p_{o}$ is homogeneous of degree 0 for large $|\xi|$.

Then take $\phi \epsilon C^{\infty}(\pi(U))$ with $\phi=1$ near $x_{o}$ and put $p(x, \xi)=p_{o}(\xi)$ $\phi(x)$. This will do the job.

The following theorem and its corollary are refinements of Corollary 4.11 (pseudo local property of $\psi D O$ ) and Corollary 4.45 (elliptic regularity theorem).

Theorem 4.52. If $P \epsilon \psi^{m}(\Omega)$ and $u \epsilon D^{\prime}(\Omega)$, then $W F(P u) \subset W F(u) \cap$ Ess. supp $P$.

Proof. If $\left(x_{o}, \xi_{o}\right) \notin$ Ess. $\operatorname{supp} P$, by Lemma 4.51, we can find a $Q \epsilon \psi^{0}$ such that $\left(x_{o}, \xi_{o}\right) \notin$ char $Q$ and Ess. supp $P \cap$ Ess. supp $Q=\phi$ Then $Q P \epsilon \psi^{-\infty}$ by Proposition 4.50, so that $Q P \epsilon C^{\infty}$. But this implies that $\left(x_{o}, \xi_{o}\right) \notin W F(P u)$.

Suppose now that $\left(x_{o}, \xi_{o}\right) \notin W F(u)$. Then there exists $A \epsilon \psi^{0}$ with $\operatorname{Au} \epsilon C^{\infty}$ and $\left(x_{o}, \xi_{o}\right) \notin$ char $A$.

115 Claim. There exist operators $B, C \epsilon \psi^{0}$ with $\left(x_{o}, \xi_{o}\right) \notin$ char $B$ and $B P=$ $C A \bmod \psi^{-\infty}$.

Granted this, $B P u=C A u+\left(C^{\infty}\right.$ function) $\epsilon C^{\infty}$ which implies that $\left(x_{o}, \xi_{o}\right) \notin W F(P u)$.

To prove the claim, let $A_{o}$ be elliptic with $\sigma_{A}=\sigma_{A}$ on a conic neighbourhood $U$ of $\left(x_{o}, \xi_{o}\right)$. Then Ess. $\operatorname{supp}\left(A-A_{o}\right) \cap U=\phi$. By Lemma 4.51, choose $B$ with $\left(x_{o}, \xi_{o}\right) \notin$ char $B$ and Ess. supp $B \subset U$. Let $E_{o}$ be a parametrix for $A_{o}$ and set $C=B P E_{o}$. Then $C A=B P E_{o} A=$ $B P E_{o}\left(A-A_{o}\right)+B P E_{o} A_{o}$. Since $E_{o} A_{o}=I \bmod \psi^{-\infty}$ and $B P E_{o}\left(A-A_{o}\right)$ also belongs to $\psi^{-\infty}$ (by Proposition 4.50 again), $C A=B P \bmod \psi^{-\infty}$. This completes the proof.

Corollary 4.53. If $P$ is elliptic, then $W F(u)=W F(P u))$.
Proof. By the theorem, $W F(P u) \subset W F(u)$. If $E$ is a parametrix for

$$
\begin{aligned}
P, \text { then } W F(u) & =W F\left(E P u+C^{\infty}\right. \text { function) } \\
& =W F(E P u) \subset W F(P u), \text { by the theorem again }
\end{aligned}
$$

Hence $\quad W F(u)=W F(P u)$.
Theorem 4.56. Let $u \in \mathcal{D}^{\prime}(D)$. Then $\left(x_{o}, \xi_{o}\right) \notin W F(u)$ if and only if there exists $\phi \epsilon C_{o}^{\infty}\left(\mathbb{R}^{n}\right)$ such that $\phi\left(x_{o}\right)=1$ and $(\phi u)^{\hat{1}}$ is rapidly decreasing on a conic neighbourhood of $\xi_{o}$.

Proof. (Sufficiency) If such a $\phi$ exists, choose $p \epsilon S^{0}, p(x, \xi)=p(\xi)$ with $p(\xi)=1$ near $\xi_{o}$ and $p(\xi)=0$ outside the region where $(\phi u)$ is rapidly decreasing. Then $p(\phi u) \epsilon S$ and hence $p(D)(\phi u) \epsilon S$ where $p(D)=$ $p(x, D)$. The operator $P u=\phi p(D)(\phi u)$ is a pseudo differential operator of order 0 with symbol $\phi(x)^{2} p(\xi)$ modulo $S^{-1}$, so $\left(x_{o}, \xi_{o}\right) \notin \operatorname{char} P$. This implies that $\left(x_{o}, \xi_{o}\right) \notin W F(u)$.
(Necessity) Suppose $\left(x_{o}, \xi_{o}\right) \notin W F(u)$. Then there exists a neighbourhood $U$ of $x_{o}$ in $\Omega$ such that $\left(x_{o}, \xi_{o}\right) \notin W F(u)$ for all $x \in U$. Choose $\phi \in C_{o}^{\infty}(U)$ such that $\phi\left(x_{o}\right)=1$. Then $\left(x_{o}, \xi_{o}\right) \notin W F(\phi u)$ for all $x \in \mathbb{R}^{n}$. Let

$$
\sum=\{\xi:(x, \xi) \epsilon W F(\phi u) \text { for some } x\} .
$$

This $\sum$ does not contain $\xi_{o}$ and is a closed conic set. There exists $p(\xi) \epsilon S^{0}, p=1$ on a conic neighbourhood of $\xi_{o}$ and $p=0$ on a neighbourhood of $\sum$, say $\sum_{o}$. Since $p(\xi)=0$ on $\sum_{o}$, Ess. supp $p(D) \cap\left(\mathbb{R}^{n} \times\right.$ $\left.\sum_{o}\right)^{c}$ and hence Ess. $\operatorname{supp} p(D) \cap W F(\phi u)=\phi$. But this gives Ess. $\operatorname{supp}(p(D) \phi u)=\phi$. So $p(D) \phi u \epsilon C^{\infty}$. We now claim that $p(D) \phi u \epsilon S$.

Accepting this, we see that $p(\phi u) \hat{\epsilon} S S$ and in particular, $(\phi u)^{\wedge}$ is rapidly decreasing near $\xi_{o}$. To prove the claim, observe that

$$
D^{\beta}\left(\xi^{\alpha} p(\xi)\right)=0(1+|\xi|)^{|\alpha|-|\beta|} \epsilon L^{2} \text { if }|\alpha|-|\beta|<-\frac{n}{2}
$$

So if we put $K(x)=\tilde{p}(x), x^{\beta} D^{\alpha} K(x) \epsilon L^{2}$ if $|\alpha|-|\beta|<-\frac{n}{2}$. An application of Leibniz's rule and the the Sobolev imbedding theorem shows that $x^{\beta} D^{\alpha} K(x) \epsilon L^{\infty}$ if $|\beta|-|\alpha|>-\frac{3 n}{2}+2$. Hence $K$ and all its derivatives are rapidly decreasing at $\infty$ and since $\phi u \epsilon E^{\prime}$, the same is true of $p(D) \phi u=$ $\phi u * K$.

## 10 Some Further Applications of Pseudo Defferential Operators

We conclude our discussion of pseudo differential operators by giving117 brief and informal descriptions of some further applications. We recall that an $m^{t h}$ order ordinary differential equation $u^{(m)}=F\left(x, u, u^{(1)}, \ldots\right.$ $u^{(m-1)}$ ) can be reduced to the first order system

$$
\frac{d}{d x}\left(\begin{array}{c}
u_{1} \\
u_{2} \\
\vdots \\
u_{m}
\end{array}\right)=\left(\begin{array}{c}
u_{2} \\
u_{3} \\
\vdots \\
F\left(x, u_{1}, \ldots, u_{m}\right)
\end{array}\right)
$$

simply by introducing the derivatives of $u$ of order $<m$ as new variables, and this reduction is frequently a useful technical device. To do
something similar for $m^{\text {th }}$ order partial differential equations, however, is more problematical. If one simply introduces all partial derivatives of $u$ of order $<m$ as new variables and writes down the first order differential relations they satisfy, one usually obtains more equations than there are unknowns, because of the equality of the mixed partials. Moreover, such a reduction usually does not preserve the character of the original equation. For example take the Laplace equation in two variables :

$$
\frac{\partial^{2} u}{\partial x^{2}}+\frac{\partial^{2} u}{\partial y^{2}}=f .
$$

putting $u_{1}=u, u_{2}=\frac{\partial u}{\partial x}, u_{3}=\frac{\partial}{\partial y}$ we have a $3 \times 3$ system given by

$$
\begin{gathered}
\frac{\partial u_{1}}{\partial x}-u_{2}=0 \\
\frac{\partial u_{1}}{\partial y}-u_{3}=0 \\
\frac{\partial u_{2}}{\partial x}+\frac{\partial u_{3}}{\partial y}=f
\end{gathered}
$$

The original equation is elliptic. But consider the matrix of the top order symbols of the $3 \times 3$ system obtained above. The matrix is given by

$$
2 \pi i\left(\begin{array}{lll}
\xi & 0 & 0 \\
\eta & 0 & 0 \\
0 & \xi & \eta
\end{array}\right)
$$

which is not invertible. This means that the first order system is not elliptic.

There is, however, a method, due to A.P. CALDERON, of reducing an $m^{t h}$ order linear partial differential equation to an $m \times m$ system of first order pseudo differential equations which preserves the characteristic variety of the equation in a sense which we shall make precise below. In this method, one of the variables is singled out to play a special role, so we shall suppose that we are working on $\mathbb{R}^{n+1}$ with coordinates $\left(x_{1}, x_{2}, \ldots, x_{n}, t\right)$.

Let $L$ be a partial differential operator of order $m$ on $\mathbb{R}^{n+1}$ such that the coefficient of $\partial_{t}^{m}$ is nowhere vanishing. Dividing throughout by this coefficient, we can assume that $L$ is of the following form:

$$
L=\partial_{t}^{m}-\sum_{j=0}^{m-1} A_{m-j}\left(x, t, D_{x}\right) \partial_{t}^{j}
$$

Here $A_{m-j}$ is a differential operator of order $\leq m-j$ in the $x$ variable with coefficients depending on $t$.

We want to reduce the equation $L v=f$ to a first order system. To this end we proceed as follows:

Using the familiar operators $\Lambda^{s}$ with symbol $(1+|\xi|)^{s / 2}$ (acting in the $x$ variables ) we put

$$
\begin{aligned}
u_{1} & =\Lambda^{m-1} v \\
u_{2} & =\Lambda^{m-2} \partial_{t^{v}} \\
u_{3} & =\Lambda^{m-3} \partial_{t}^{2} v \\
\vdots & \\
u_{m} & =\partial_{t}^{m-1} v .
\end{aligned}
$$

For $j<m$, we observe that $\partial_{t} u_{j}=\Lambda u_{j+1}$. Therefore,

$$
\begin{aligned}
\partial_{t} u_{m} & =f+\sum_{j=0}^{m-1} A_{m-j}\left(x, t, D_{x}\right) \partial_{t}^{j} v \\
& =f+\sum_{j=0}^{m-1} A_{m-j}\left(x, t, D_{x}\right) \Lambda^{j-m+1} u_{j+1} \\
& =f+\sum_{j=1}^{m} A_{m-j}\left(x, t, D_{x}\right) \Lambda^{j-m} u_{j}
\end{aligned}
$$

Set $B_{j}\left(x, t, D_{x}\right)=A_{m-j+1}\left(x, t, D_{x}\right) \Lambda^{j-m}$. This $B_{j}$ is a $\Psi D O$ of order 1 in the variable $x$. Then the equation $L v=f$ is equivalent to

$$
\partial_{t} u=K u+\tilde{f}
$$

where $u=\left(u_{1}, \ldots, u_{m}\right)^{t}, \tilde{f}=(0,0, \ldots, f)^{t}$. (Here $(\cdots)^{t}$ denotes the transpose of the row vector $(\cdots)$ ) and $K$ is the matrix given by

$$
K=\left(\begin{array}{ccc}
0 & \Lambda & 0 \cdots 0 \\
0 & 0 & \Lambda \cdots 0 \\
. & . & \cdots \\
. & . & \cdots \\
0 & 0 & 0 \cdots \Lambda \\
B_{1} & B_{2} & B_{3} \cdots B_{m}
\end{array}\right)
$$

This $K$ is a matrix of first order pseudo differential operators in $x$, with coefficients depending on $t$.

Let us now make a little digression to fix on some notations which will be used in the further development. For a $p \epsilon S^{m}$, a function $p_{m}(x, \xi)$ homogeneous of degree $m$ in $\xi$ is said to be the Principal symbol of $p(x, D)$ if $p-p_{m}$ agrees with an element of $S^{m-1}$ for large $|\xi|$. We remark that not all $\Psi D O^{\prime} s$ have a principal symbol; but most of them that arise in practice do. Moreover, the principal symbol, if it exists, is clearly unique.

## Examples

(i) If $p$ is a polynomial,

$$
p(x, \xi)=\sum_{|\alpha| \leq m} a_{\alpha}(x) \xi^{\alpha}
$$

then

$$
p_{m}(x, \xi)=\sum_{|\alpha|=m} a_{\alpha}(x) \xi^{\alpha}
$$

(ii) If $p(x, \xi)=\left(1+|\xi|^{2}\right)^{s / 2}, p_{s}(x, \xi)=|\xi|^{s}$.

Returning to our discussion, let $a_{k}(x, t, \xi)$ be the principal symbol of $A_{k}\left(x, t, d_{x}\right)$ (Here we are considering $A_{k}$ as an operator of order $k$, so if it happens to be of lower order, its principal symbol is zero). Then
the principal symbol of $B_{j}$ is $b_{j}(x, t \xi)=a_{m-j+1}(x, t, \xi)|\xi|^{j-m}$ and the principal symbol of $k$ is the matrix $K_{1}$

$$
K_{1}(x, t, \xi)=\left(\begin{array}{cccccc}
0 & |\xi| & 0 & 0 & \cdots & 0 \\
0 & 0 & |\xi| & 0 & \cdots & 0 \\
0 & 0 & 0 & |\xi| & \cdots & 0 \\
. & . & . & & \cdots & \\
. & . & . & & \cdots & \\
. & . & . & & \cdots & \\
0 & 0 & 0 & 0 & \cdots & |\xi| \\
b_{1} & b_{2} & b_{3} & b_{4} & \cdots & b_{m}
\end{array}\right)
$$

The principal symbol of $L$, on the other hand, is

$$
L_{m}(x, t, \xi, \tau)=(2 \pi i \tau)^{m}-\sum_{j=0}^{m-1} a_{m-j}(x, t, \xi)(2 \pi i \tau)^{j}
$$

The characteristic variety of $L$, i.e., the set of zeros of $L_{m}$, can easily be read off from the matrix $K$ as follows:

Proposition 4.57. The eigenvalues of $K_{1}(x, t, \xi)$ are precisely $2 \pi i \tau_{1}, \ldots$, $2 \pi i \tau_{m}$ where $\tau_{1}, \tau_{2}, \ldots, \tau_{m}$ are the roots of the polynomial $L_{m}(x, t, \xi,$.$) .$ Proof. The characteristic polynomial $p(\lambda)$ of $K_{1}$ is the determinant of

$$
\left(\begin{array}{ccccc}
\lambda & -|\xi| & 0 & \cdots & 0 \\
0 & \lambda & -|\xi| & \cdots & 0 \\
\vdots & & & & \\
0 & 0 & 0 & \cdots & -|\xi| \\
-b_{1} & -b_{2} & -b_{3} & : & \lambda-b_{m}
\end{array}\right)
$$

Expanding in minors along the last row, we get

$$
\begin{aligned}
p(\lambda) & =\sum_{j=1}^{m-1}(-1)^{m+j}\left(-b_{j}\right) \lambda^{j-1}(-|\xi|)^{m-j}+\lambda^{m-1}\left(\lambda-b_{m}\right) \\
& =\lambda^{m}-\lambda^{m-1} b_{m}-\sum_{j=1}^{m-1} b_{j}|\xi|^{m-1} \lambda^{j-1}
\end{aligned}
$$

$$
\begin{aligned}
& =\lambda^{m}-\sum_{j=1}^{m} a_{m-j+1} \lambda^{j-1}=\lambda^{m}-\sum_{j=0}^{m-1} a_{m-j} \lambda^{j} \\
& =L_{m}\left(x, t, \xi,(2 \pi i)^{-1} \lambda \lambda .\right.
\end{aligned}
$$

This proves the proposition.
It is also easy to incorporate boundary conditions into this scheme. Suppose we want to solve the equation $L v=f$ with the boundary conditions $B_{j} v_{t=0}=g_{j}, 1,2, \ldots, v$ where

$$
B_{j}=\sum_{k=0}^{m_{j}} b_{j}^{k}\left(x, t, D_{x}\right) \partial_{t}^{k}, b_{j}^{k} \text { is of order } m_{j}-k \text { and } m_{j} \leq m-1 .
$$

If we set

$$
\begin{aligned}
B_{j}^{k} & =\Lambda^{m-m_{j}-1} b_{j}^{k-1}\left(x, 0, D_{x}\right) \Lambda^{k-m}, \\
\phi_{j} & =\Lambda^{m-m_{j}-1} g_{j}
\end{aligned}
$$

then with $u_{j}=\Lambda^{m-j} \partial_{t}^{j-1} v$ as above and $u_{j}^{0}(x)=u_{j}(x, 0)$, the boundary conditions will become

$$
\sum_{k=1}^{m_{j}+1} B_{j}^{k} u_{k}^{0}=\phi_{j}, j=1,2, \ldots v .
$$

This is a system of zeroth order pseudo differential equations.
The above method of reduction is useful, for example, in the following two important problems.

## (1) Cauchy Problem for Hyperbolic Equations

$$
L v=f,\left.\partial_{t}^{j} v\right|_{t=0}=g_{j}, j=0,1, \ldots, m-1 .
$$

Here the equation is said to be hyperbolic when the eigenvalues of the matrix $K_{1}$ (i.e., the principal symbol of $K$ occurring in the linear system corresponding to $L v=f$ ) are purely imaginary.Equivalently, the roots of $L_{m}(x, t, \xi,$.$) are real.$

## (2) Elliptic Boundary Value Prob Lems

$$
L u=f \text { on } \Omega, B_{j} u=g_{j} \text { on } \partial \Omega
$$

where $L$ is elliptic of order $2 m$ and $j=1,2, \ldots m$. Here one works locally near a point $x_{0} \epsilon \partial \Omega$ and makes a change of coordinates so that $\partial \Omega$ becomes a hyperplane near $x_{0}$. One then can apply Calderon's reduction technique, taking $t$ as the variable normal to $\partial \Omega$. (See Michael Taylor [3]).

We shall now sketch an example of a somewhat different technique for applying $\psi D O^{\prime} S$ to elliptic boundary value problems.

Let $\Omega$ be a bounded open set of $\mathbb{R}^{n}$ with $C^{\infty}$ boundary $\partial \Omega$ Consider the problem

$$
\Delta u=0 \text { in } \Omega, \partial_{X} u+a u=g \text { on } \partial \Omega, a \epsilon C^{\infty}(\partial \Omega) .
$$

Here $X$ is a real a vector field on the boundary which is nowhere tangent to the boundary and $\partial_{X} u=\operatorname{grad} u \cdot \chi$.

If $v$ is the unit outward normal to $\partial \Omega$, by normalising we can assume that $\chi=v+\tau$ where $\tau$ is tangent to $\partial \Omega$. We want to use $\Psi D O$ to reduce this to the Dirichlet problem. Pretend for the moment that $\partial \Omega=$ $\mathbb{R}^{n-1} x\{0\}$ and $\Omega=\left\{x: x_{n}<0\right\}$ so that $\partial_{v}=\frac{\partial}{\partial x_{n}}$.

If $\Delta u=0$ then $\partial_{v}^{2} u=-\sum_{j=1}^{n-1} \frac{\partial^{2} u}{\partial x_{j}^{2}}=-\Delta_{b} u, \Delta_{b}$ being the Laplacian on the boundary. So formally $\partial_{\nu} u= \pm \sqrt{ }\left(\Delta_{b} u\right)$. Indeed, if we compare this with our discussion of the Poisson kernel in Chapter 2(taking account of the fact that $\Omega$ is now the lower rather than the upper half space), we see that the equation $\partial_{v} u=\sqrt{ }\left(-\Delta_{b}\right) u$ is correct if interpreted in terms of the Fourier transform in the variables $x^{\prime}=\left(x_{1}, \ldots x_{n-1}\right)$;

$$
\partial_{\nu} \tilde{u}\left(\xi^{\prime}, x_{n}\right)=2 \pi\left|\xi^{\prime}\right| \tilde{u}\left(\xi^{\prime}, x_{n}\right) .
$$

It turns out that something similar works for our original domain $\Omega$. Namely, if $u$ is smooth on $\bar{\Omega}$ and $\Delta u=0$ on $\Omega$ then $\left.\partial_{\nu} u\right|_{\partial \Omega}=P\left(\left.u\right|_{\partial \Omega)}\right.$ where $P$ is a pseudo differential operator of order 1 on $\partial \Omega$ which equals $\sqrt{ }\left(-\Delta_{b}\right)$ modulo terms of order $\leq 0$ where $\Delta_{b}$ is the Laplace -Beltrami operator on $\partial \Omega$. In particular, $P$ is elliptic and has real principal symbol.

The boundary conditions become $P v+\left|\partial_{t} v+a v=g, v=u\right|_{\partial \Omega}$. This is a first order pseudo differential equation for $v$. It is elliptic (because $\partial_{\tau}$ has imaginary symbol). So it can be solved modulo smoothing operators. Since $\partial \Omega$ is compact. smoothing operators on $\partial \Omega$ are compact, so our boundary equation can be solved provided $g$ is orthogonal to some finite dimensional space of smooth functions.

Having done this, we have reduced our original problem to the familiar Dirichlet problem $\Delta u=0$ in $\Omega$ and $\left.u\right|_{\partial \Omega}=v$ which has a unique solution.

## Chapter 5

## $L^{P}$ and Lipschitz Estimates

OUR AIM IN this chapter is to study how to measure the smoothing properties of pseudo differential operators of non positive order in terms of various important function spaces. Most of the interesting results are obtained by considering operators of order $-\lambda$ with $0 \leq \lambda \leq n$. Indeed, if $P \epsilon \Psi^{-\lambda}$ with $\lambda>n$ and $K(x, y)$ is the distribution kernel of $P$ then, by Theorem 4.10 we know that $K \epsilon C^{j} \Omega \times \Omega$ ) when $j<\lambda-$ $n$. So these operators can be studied by elementary methods. What is more, when $P \epsilon \Psi^{-\lambda}, D^{\alpha} P \epsilon \Psi^{-\lambda+|\alpha|}$. So, by a proper choice of $\alpha$, we can make $0 \leq \lambda-|\alpha| \leq n$ and then study $D^{\alpha} P$ rather than $P$ itself. Actually, we shall restrict attention to operators of order $-\lambda$ where $0 \leq$ $\lambda<n$. The transitional case $\lambda=n$ requires a separate treatment to obtain sharp results; however, for many purpose, it suffices to make the trite observation that an operator of order $-n$ can be regarded as an operator of order $-n+\epsilon$. Further, we restrict our attention to $\Psi D O^{\prime} s$ whose symbols have asymptotic expansions

$$
p(x, \xi) \sim \sum_{j=0}^{\infty} p_{j}(x, \xi)
$$

where $P_{j}$ is homogeneous of degree $-\lambda-j$. (These $p_{j}^{\prime} s$ no longer belong to our symbol classes, being singular at $\xi=0$, but we can still consider the operators $p_{j}(x, D)$.) By the preceding remarks, it will suffice to consider the operators corresponding to the individual terms in the series
$\sum p_{j}$ whose degrees of homogeneity are between 0 and $-n$. Thus we are looking at $p(x, D)$ where $p(x, \xi)$ is $C^{\infty}$ on $\mathbb{R}^{n} \times \mathbb{R}^{n} \backslash\{0\}$ and homogeneous of degree $-\lambda$ in $\xi$ where $0 \leq \lambda<n$.

Since $\lambda<n$, for each $x, p(x$, .) is locally integrable at the origin and hence defines a tempered distribution. Denoting the inverse Fourier transform of this distribution $p_{2}^{\vee}(x,$.$) we then have$

$$
\begin{aligned}
p(x, D) u(x) & =\int e^{2 \pi i x . \xi} p(x, \xi) \hat{u}(\xi) d \xi \\
& =\left(p_{2}^{\vee}(x, .) * u\right)(x) \\
& =\int p_{2}^{\vee}(x, x-y) u(y) d y .
\end{aligned}
$$

Thus we see that the distribution kernel of $p(x, D)$ is given by $K(x, y)$ $=p_{2}^{v}(x, x-y)$.

Let us digress a little to make some remarks on homogeneous distributions.

Definition 5.1. A distribution $f \in S^{\prime}$ is said to be homogeneous of degree $\mu$, if $\left\langle f, \phi_{r}\right\rangle=r^{\mu}<f, \phi>$ for all $\phi \epsilon S$ where $\phi_{r}$ is the function defined by $\phi_{r}(x)=r^{-n} \phi(x / r)$.

## Exercise.

1. Show that the above definition agrees with the usual definition of homogeneity when $f$ is a locally integrable function.
2. Show that if $f \epsilon S^{\prime}$ is homogeneous of degree $\mu$, then $D^{\alpha} f$ is homogeneous of degree $\mu-|\alpha|$.
3. Show that $D^{\alpha} \delta$ is homogeneous of degree $-n-|\alpha|$ (where $\delta$ is the Dirac measure at 0 ).

Let us now prove a proposition concerning the Fourier transform of homogeneous distributions. It turns out that the Fourier transform of such a distribution is also a homogeneous one. Precisely, we have the following

Proposition 5.2. If $f$ is a tempered distribution, homogeneous of degree $\mu$, then $\hat{f}$ is homogeneous of degree $-\mu-n$. If $f$ is also $C^{\infty}$ away from the origin, then the same is true of $\hat{f}$.
Proof If $f$ is homogeneous of degree $\mu$, then for $\phi \epsilon S$

$$
\begin{aligned}
<\hat{f}, \phi_{r}> & =<f,\left(\phi_{r}\right)^{\wedge}>=<f, r^{-n} \hat{\phi}_{1 / r}>=r^{-n-\mu}<f, \hat{\phi}> \\
& =r^{-\mu-n}<\hat{f}, \phi>
\end{aligned}
$$

This proves the first assertion. To prove the second assertion, choose $\phi \epsilon C_{0}^{\infty}$ with $\phi=1$ in a neighbourhood of the origin and write $f=\phi f+$ $(1-\phi) f$. Since $\phi f \in E^{\prime},(\phi f)^{)}$is $C^{\infty}$ everywhere. On the other hand, $(1-\phi) f$ is $C^{\infty}$ and homogeneous of degree $\mu$ for large $\xi$, and hence lies in $S^{\mu}\left(\mathbb{R}^{n}\right)$. Let $p(x, \xi)=(1-\phi)(\xi) f(\xi)$. The distribution kernel of $p(x, D)$ is given by

$$
\begin{aligned}
K(x, y) & =\int e^{-2 \pi i(y-x), \xi}(1-\phi)(\xi) f(\xi) d \xi \\
& =((1-\phi) f)(y-x) .
\end{aligned}
$$

Since $K$ is $C^{\infty}$ away from the diagonal, we see that $((1-\phi) f)$ is $C^{\infty}$ away from the origin. Hence $\hat{f}$ is $C^{\infty}$ on $\mathbb{R}^{n} \backslash\{0\}$.

Returning to our operators $p(x, D)$ with $p(x, \xi)$ homogeneous of de-
gree $-\lambda$ in $\xi, 0 \leq \lambda<n$, by the proposition above, we have

$$
p(x, D) u=\int K(x, x-y) u(y) d y
$$

where $K=P_{2}^{\vee}$ is $C^{\infty}$ on $\mathbb{R}^{n} \times \mathbb{R}^{n} \backslash\{0\}$ and homogeneous of degree $\lambda-n$ in the second variable.

Finally, we restrict attention to constant coefficient case, i.e., $K(x, x-$ $y)=K(x-y)$. The essential ideas are already present in this case and the results we shall obtain can be generalised to the variable coefficient case in a rather routine fashion. Let us give a name to the objects we are finally going to study.
Definition 5.3. A tempered distribution $K$ which is homogeneous of degree $\lambda-n$ and $C^{\infty}$ away from the origin is called a kernel of type $\lambda$. If $K$ is a kernel of type $\lambda$, the operator $T f=K * f$ is called an operator of type $\lambda$.

We now classify the kernels of type $\lambda \geq 0$.
Proposition 5.4. Suppose $\lambda>0$ and $f \epsilon C^{\infty}\left(\mathbb{R}^{n} \backslash\{0\}\right)$ is homogeneous of degree $\lambda-n$ in the sense of functions. Then $f$ is locally integrable and defines a distribution $F$ which is a kernel of type $\lambda$.

Conversely, if $F$ is a kernel of type $\lambda$ and $f$ is the function with which $F$ agrees on $\mathbb{R}^{n} \backslash\{0\}$, then $\langle F, \phi\rangle=\int f \phi$ for every $\phi$ in $S$.

Proof. Since $\lambda>0, f$ is locally integrable and of (at most ) polynomial growth at $\infty$, so $f$ defines an $F$ in $S^{\prime}$. It is easy to check that $F$ is
130 homogeneous in the sense defined above, i.e. $\left.\left.<F, \phi_{r}\right\rangle=r^{\lambda-n}<F, \phi\right\rangle$, and so $F$ is a kernel of type $\lambda$.

For the converse, define $G$ by $<G, \phi>=<F, \phi>-\int f \phi, \phi \in S$. Then $G$ is a distribution supported at 0 and hence we have $G=\sum c_{\alpha} D^{\alpha} \delta$. Therefore,

$$
<G, \phi_{r}>=\sum c_{\alpha} r^{-n-|\alpha|}\left(D^{\alpha} \phi\right)(0)=O\left(r^{-n}\right), \text { as } r \rightarrow \infty
$$

On the other hand,

$$
<G, \phi_{r}>=r^{-n-\lambda}<G, \phi>
$$

and this is not $O\left(r^{-n}\right)$ as $r$ tends to $\infty$ unless $\langle G, \phi\rangle=0$. Hence $G=0$ and this completes the proof.

Suppose that $f \epsilon C^{\infty}\left(\mathbb{R}^{n} \backslash\{0\}\right)$ is homogeneous of degree $-n$. Then $f$ is not locally integrable near 0 and so does not define a distribution in a trivial way. However, let us define

$$
\mu_{f}=\int_{|x|=1} f(x) d \sigma(x)
$$

If $\mu_{f}=0$ there is a canonical distribution associated with $f$ which is called principal value of $f, P V(f)$, defined by

$$
<P V(f), \phi>=\lim _{\epsilon \rightarrow 0} \int_{|x|>\epsilon} f(X) d x
$$

To see that this limit exists, we observe that

$$
\int_{\epsilon<|x|<1} f(x) d x=\mu_{f} \int_{\epsilon}^{1} r^{-1} d r=0
$$

Hence

$$
<P V(f), \phi>=\lim _{\epsilon \rightarrow 0} \int_{\epsilon<|x|<1} f(x)(\phi(x)-\phi(0)) d x+\int_{|x| \geq 1} f(x) \phi(x) d x
$$

where the last integrals are absolutely convergent, since

$$
\begin{gathered}
|\phi(x)-\phi(0)| \leq c|x|, \quad \text { so that } \\
\int_{|x|<1}|f(x) \| \phi(x)-\phi(0)| d x \leq c \int_{|x|<1}|x|^{-n+1} d x<\infty
\end{gathered}
$$

Further, the estimate on $\phi(x)-\phi(0)$ depends only on the first derivatives of $\phi$ via the mean value theorem, so it is easily verified that the functional $P V(f)$ is continuous on $S$. Finally, we observe that

$$
\begin{aligned}
<P V(f), \phi_{r}> & =\lim _{\epsilon \rightarrow 0} \int_{|x|>\epsilon} f(x) r^{-n} \phi(x / r) d x \\
& =\lim _{\epsilon \rightarrow 0} \int_{|y|>(\epsilon / r)} f(y) r^{-n} \phi(y) d y \\
& =r^{-n}<P V(f), \phi>
\end{aligned}
$$

Thus $P V(f)$ is homogeneous of degree $-n$ and so is a kernel of type 0.

The following theorem gives a sort of converse to the above results.
Theorem 5.5. Suppose $F$ is a kernel of type 0 and $f$ is the function with which $F$ agrees on $\mathbb{R}^{n} \backslash\{0\}$. Then $\mu_{f}=0$ and $F=P V(f)+c \delta$, for some constant $c$.

Proof. Define the functional $G$ on $S$ by

$$
<G, \phi>=\int_{|x| \leq 1} f(x)\left(\phi(x)-(\phi(0)) d x+\int_{|x|>1} f(x) \phi(x) d x\right.
$$

By the same argument as above, $G$ is a tempered distribution and $G=F$ on $\mathbb{R}^{n} /\{0\}$. Therefore, $G-F=\sum c_{\alpha} D^{\alpha} \delta$. Now, $<F, \phi_{r}>=r^{-n}<$ $P, \phi>$ so that

$$
<G, \phi_{r}>-r^{-n}<G, \phi>=<G-F \phi_{r}>-r^{-n}<G-F, \phi>=0\left(r^{-n}\right)
$$

as $r$ tends to $\infty$. On the other hand, for $r>1$,

$$
\begin{aligned}
<G, \phi_{r}> & -r^{-n}<G, \phi>=\int_{|x| \leq 1 / r} r^{-n} f(x)(\phi(x)-\phi(0)) d x \\
& +\int_{|x|>1 / r} r^{-n} f(x) \phi(x) d x-r^{-n} \int_{|x| \leq 1} f(x)(\phi(x)-\phi(0)) d x-r^{-n} \\
& \int_{|x|>1} f(x) \phi(x) d x \\
& =r^{-n} \phi(0) \int_{1 r \leq|x| \leq 1} f(x) d x=r^{-n} \phi(0) \log r \mu_{f^{\prime}} \text { for every } \phi
\end{aligned}
$$

This is not $0\left(r^{-n}\right)$ as $r$ tends to $\infty$ unless $\mu_{f}=0$. Finally, $F-P V(f)$ is a kernel of type 0 which is supported at the origin and hence is a multiple of $\delta$.

To study the boundedness of operators of type $\lambda$ on $L^{p}$ spaces, we need some concepts from measure theory.
Definition 5.6. Let $F$ be a measurable function defined on $\mathbb{R}^{n}$. Then the distribution function of $f$ is the function $\delta_{f}:(0, \infty) \rightarrow[0, \infty]$ given by $\delta_{f}(t)=\left|E_{t}\right|$, where $E_{t}=\{x:|f(x)|>t\}$.

From the distribution function of $f$, we can get a large amount of information regarding $f$. For example, we have

$$
\int|f(x)|^{p} d x=-\int_{0}^{\infty} t^{p} d \delta_{f}(t)
$$

(To see this, observe that the Riemann sums for the Stieltjes integral on the right are approximating sums for the Lebesque integral on the left).

If $t^{p} \delta_{f}(t)$ converges to 0 as $t$ tends to 0 and $t$ tends to $\infty$, we can integrate by parts to obtain $\int|f(x)|^{p} d x=p \int_{0}^{\infty} t^{p-1} \delta_{f}(t) d t$.

Using the concept of distribution functions, we will now define weak $\mathbf{1 3 3}$ $L^{p}$ spaces.

Definition 5.7. For $1 \leq p<\infty$ we define weak $L^{p}$ as $\left\{f: \delta_{f}(t) \leq(c / t)^{p}\right.$ for some constant $c\}$. For $f \in$ weak $L^{p}$, the smallest such constant $c$ will be denoted by $[f]_{p}$. Thus, if feweak $L^{p}$, we have

$$
\delta_{f}(t) \leq\left([f]_{p} / t\right)^{p}
$$

For $p=\infty$, we set weak $L^{\infty}=L^{\infty}$.
Proposition 5.8. CHEBYSHEV 'S INEQUALITY) $L^{p} \subset$ weak $L^{p}$ and $[f]_{p} \leq\|f\|_{p}$.

Proof. For $f \epsilon L^{p}$, if $E_{t}=\{x:|f(x)|>t\}$,

$$
\|f\|_{p}^{p}=\int_{\mathbb{R}^{n}}|f(x)|^{p} d x \geq \int_{E_{t}}|f(x)|^{p} d x \geq t^{p}\left|E_{t}\right|
$$

i.e., $\quad\left|E_{t}\right| \leq\left(\|f\|^{p} / t\right) p$.

From this, it follows that $f \in$ weak $L^{p}$ and $[f]_{p} \leq\|f\|_{p}$.
Remark 5.9. It is not true that $L^{p}=$ weak $L^{p}$ for $1 \leq p<\infty$. For example, $f(x)=|x|^{-n / p}$ belongs to weak $L^{p}$ but not $L^{p}$.

Remark 5.10. The function $f \rightarrow[f]_{p}$ satisfies $[c f]_{p}=|c|[f]_{p}$.
But it fails to satisfy triangle inequality and hence is not a norm. However, since

$$
\{|f+g|>t\} \subset\{|f| t / 2\} \cup\{|g|>t / 2\}
$$

we have

$$
\delta_{f+g}(t) \leq \delta_{f}(t / 2)+\delta_{g}(t / 2)
$$

which gives, when $f$ and $g$ are in weak $L^{p}$,

$$
\begin{gathered}
\delta_{f+g}(t) \leq\left(2^{p}[f]_{p}^{p}+2^{p}[g]_{p}^{p}\right) / t^{p} \text { so that } \\
{[f+g]_{p} \leq 2\left([f f]_{p}^{p}+[g]_{p}^{p}\right)^{1 / p} \leq 2\left([f f]^{p}+[g]_{p}\right) .}
\end{gathered}
$$

The functional $[.]_{p}$ thus defines a topology on weak $L^{p}$ which will turn it into a (non-locally convex) topological vector space.

Definition 5.11. A linear operator $T$ defined on a space of functions is said to be of weak type $(p, q), 1 \leq p, q \leq \infty$, if $T$ is a bounded linear operator from $L^{p}$ into weak $L^{q}$, i.e. for every $f \epsilon L^{p},[T f]_{q} \leq c\|f\|_{p}$ for some constant $c$ independent of $f$.

We will now state the Marcinkiewicz interpolation theorem and use it to prove generalisations of Young's inequality (Theorem 1.3).

Theorem 5.12 (J. Marcinkiewicz). Suppose $T$ is of weak types ( $p_{o}, q_{o}$ ) and $\left(p_{1}, q_{1}\right)$ with $1 \leq p_{i} \leq q_{i} \leq \infty, p_{0}<p_{1}, q_{0} \neq q_{1}$, i.e., $\left[{ }^{\prime} \Gamma f\right]^{q_{i}} \leq$ $c_{i}\|f\|_{p_{i}}$ for $i=0,1$. Then, if

$$
1 / p_{\theta}=(1-\theta) / p_{0}+\left(\theta / p_{1}\right), 1 / q_{\theta}=(1-\theta) / q_{0}+\left(\theta / q_{1}\right), 0<\theta<1
$$

$T$ is bounded from $L^{p_{\theta}}$ to $L^{q_{\theta}}$ i.e., $\|T f\|_{q_{\theta}} \leq c_{\theta}\|f\|_{p_{\theta}}$ where the constant $c_{\theta}$ depends only on $p_{0}, q_{0}, p_{1}, q_{1}, c_{0}, c_{1}$ and $\theta$.

For the proof of this theorem, see A. Zygmund [5] or E.M.Stein [2].
Theorem 5.13 (GENERAL FROM OF YOUNG'S INEQUALITY). If $(1 / p)+(1 / q)-1=1 / r, 1 \leq p, q, r<\infty$, Then, for $f \epsilon L^{p}, g \epsilon L^{q} f * g \epsilon L^{r}$ and we have $\|f * g\|_{r} \leq c_{p q}\|f\|_{p}\|g\|_{q}$.
(In fact, $c_{p q} \leq 1$ for all $p, q$, although our proof does not yield this estimate).

Proof. Fixing $f \epsilon L^{p}$, consider the convolution operator $g \rightarrow T_{q}=f * g$. We know that for $g \epsilon L^{1}, T g \epsilon L^{p}$ and $\|T g\|_{p} \leq\|f\|_{p}\|g\|_{1}$. Further if $p^{\prime}$ is the conjugate of $p$, then, by Hölder's inequality $\|T g\|_{\infty} \leq\|g\|_{p},\|f\|_{p}$ for $g \epsilon L^{p^{\prime}}$. Thus we see that the operator $T$ is of weak types $(1, p)$ and
( $p^{\prime}, \infty$ ). Therefore, by the Marcinkiewicz interpolation theorem, $T$ maps $L^{p_{\theta}}$ boundedly into $L^{q_{\theta}}$ for every $0<\theta<1$ with

$$
\left(1 / p_{\theta}\right)=1-\theta+\left(\theta / p^{\prime}\right)=1-(\theta / p) \text { and } 1 / q_{\theta}=(1 / p)-\theta / p
$$

Given $r$ with $(1 / p)+(1 / q)-1=(1 / r)$, set $q_{\theta=r}$. Then $\left(1 / p_{\theta}\right)=$ $1+(1 / r)-(1 / p)=(1 / q)$. Hence we get the required result.

Theorem 5.14 (WEAK TYPE YOUNG'S INEQUALITY). Let $1 \leq p<$ $q<\infty,(1 / p)+(1 / q)>1$ and $(1 / p)+(1 / q)-1=(1 / r)$. Suppose $f \epsilon L^{p}$ and $g \epsilon$ weak $L^{q}$. Then, we have
a) $f * g$ exists a.e. and is in weak $L^{r}$, and

$$
[f * g]_{r} \leq c_{p q}\|f\|_{p}[g]_{q}
$$

b) If $p>1$, then $f * g \epsilon L^{r}$ and $\|f * g\|_{r} \leq c_{p q}^{\prime}\|f\|_{p}[g]_{q}$.

Proof. a) It suffices to assume that $\|f\|_{p}=[g]_{p}=1$ and to show that $[f * g]_{r} \leq c_{p q}$. Given $\alpha>0$, let

$$
M=(\alpha / 2)^{p^{\prime} /\left(p^{\prime}-q\right)}\left(p^{\prime} /\left(p^{\prime}-q\right)^{-1 /\left(p^{\prime}-q\right)}\right.
$$

where $p^{\prime}$ is as usual the conjugate of $p$. Define

$$
\begin{aligned}
g_{1}(x) & =g(x), \text { if }|g(x)|>M \\
& =0 \text { otherwise } \\
\text { and } \quad g_{2}(x) & =g(x)-g_{1}(x) . \quad \text { Then } \\
\delta_{f * g}(\alpha) \leq & \delta_{f * g_{1}}(\alpha / 2)+\delta_{f * g_{2}}(\alpha / 2)
\end{aligned}
$$

and we shall estimate the quantities on the right separately. By Holder's inequality,

$$
\left\|f * g_{2}\right\|_{\infty} \leq\|f\|_{p}\left\|g_{2}\right\|_{p^{\prime}}=\left\|g_{2}\right\|_{p^{\prime}}
$$

Since $(1 / q)-\left(1 / p^{\prime}\right)=1 / r>0$, we see that $p^{\prime}-q>0$, and hence

$$
\begin{aligned}
\left\|g_{2}\right\|_{p^{\prime}}^{p^{\prime}} & =p^{\prime} \int_{0}^{\infty} t^{p^{\prime}-1} \delta_{g_{2}}(t) d t \\
& =p^{\prime} \int_{0}^{M} t^{p^{\prime}-1}\left(\delta_{g}(t)-\delta_{g}(M)\right) d t \\
& \leq p^{\prime} \int_{0}^{M} t^{p^{\prime}-1} t^{-q_{d t}}=\left(p^{\prime} /\left(p^{\prime}-q\right)\right) M^{p^{\prime}-q}=(\alpha / 2)^{p^{\prime}}
\end{aligned}
$$

Thus $\left(f * g_{2}\right)(x)$ exists at all points and $\left\|f * g_{2}\right\|_{\infty} \leq \alpha / 2$.
Consequently $\delta_{f * g_{2}}(\alpha / 2)=0$. Next consider

$$
\begin{aligned}
\left\|g_{1}\right\|_{1} & =\int_{0}^{\infty} \delta_{g_{1}}(t) d t=\int_{0}^{M} \delta_{g}(M) d t+\int_{M}^{\infty} \delta_{g}(t) d t \\
& \leq \int_{0}^{M} M^{-q} d t+\int_{M}^{\infty} t^{-q} d t
\end{aligned}
$$

The integral $\int_{M}^{\infty} t^{-q} d t$ converges since $q>1$, and we obtain

$$
\left\|g_{1}\right\|_{1} \leq M^{1-q}+M^{1-q} /(q-1)=(q /(q-1)) M^{1-q}
$$

Also, by Chebyshev inequality,

$$
\begin{aligned}
\delta_{f * g_{1}}\left(\frac{\alpha}{1}\right) & \leq\left(\frac{\left\|f * g_{1}\right\|}{\alpha / 2} p\right) p \\
& \leq\left(2^{p} / \alpha^{p}\right)(q /(q-1))^{p}(\alpha / 2)^{\frac{p^{\prime}}{p^{\prime}-q} p(1-q)}\left(p^{\prime} /\left(p^{\prime}-q\right)\right)^{-\frac{p(1-q)}{p^{\prime}-q}} \\
& =c_{p q} \alpha^{-r}
\end{aligned}
$$

Hence $(a)$ is proved.
b) The operator $f \rightarrow f * g$ is of weak type $(1, q)$ by (a). Also, if we choose $\bar{p}>p$ with $(1 / \bar{p})+(1 / q)>1$ and put $(1 \bar{r})=(1 / \bar{p})+$ $(1 / q)-1$ then by $(a), f \rightarrow f * g$ is of weak type $(\bar{p}, \bar{r})$. Therefore, by Marcinkiewicz, $f \rightarrow f * g$ is bounded from $L^{p}$ into $L^{q_{\theta}}$ with

$$
1 / p_{\theta}=1-\theta+(\theta / \bar{p}),\left(1 / q_{\theta}\right)=((1-\theta) / q)+(\theta / \bar{r}) .
$$

Putting $q_{\theta}=r$ we get $p_{\theta}=p$. Hence $T$ maps $L^{p}$ continuously to $L^{r}$ and $\|f * g\|_{r} \leq c_{p q}^{\prime}\|f\|_{p}[g]_{q}$.

Corollary 5.15. If $T$ is an operator of type $\lambda, 0<\lambda<n$, then $T$ is bounded from $L^{p}$ into $L^{q}$, whenever $1<p<q<\infty$ and $1 / q=1 / p-\lambda / n$. Also, $T$ is of weak type $(1, n /(n-\lambda)$ ).
Proof. If $K$ is a kernel of type $\lambda$, we have

$$
|K(x)|=\left.|x|^{\lambda-n} K(x /|x|)|\leq c| x\right|^{\lambda-n},
$$

which implies that $K \epsilon$ weak $L^{n /(n-\lambda)}$. Therefore, by Theorem5.14(b), if $f \in L^{p}$ and $T f=K * f$ then $T f \in L^{q}$ where
$(1 / q)=(1 / p)+\left(\left(n-\lambda(/ n)-1=(1 / p)-(\lambda / n)\right.\right.$ and $\|T f\|_{q} \leq c[f]_{p} \leq c\|f\|_{p}$.

Also we see that $T$ is of weak type $(1, n /(n-\lambda))$ by Theorem 5.14 (a).

The limiting case of this result with $\lambda=0$ is also true, but this is a $\mathbf{1 3 8}$ much deeper theorem:
Theorem 5.16. (Calderon-Zygmund) Operators of type 0 are bounded on $L^{p}, 1<p<\infty$.

Proof. Let $T$ be an operator of type 0 with kernel $K$ i.e., $T f=K * f$. To begin with, we can regard $T$ as a map from $c_{o}^{\infty}$ to $c^{\infty}$ and we shall show that $\|T f\|_{p} \leq c_{p}\|f\|_{p}$ for $f \in c_{o}^{\infty}, 1<p<\infty$, so that $T$ extends uniquely to a bounded operator on $L^{p}$.

We have $K=P V(k)+c \delta$ so that $T f=P V(k) * f+c f$. Since the identity operator is continuous, we shall assume that $c=0$ and also we shall identify $K$ with $k$. The proof now proceed in several steps.

Step 1. $T$ is bounded on $L^{2}$. Indeed, we have $(T f) t=\hat{k} \hat{f} . \hat{k}$ is smooth away from 0 and homogeneous of degree 0 , hence is bounded on $\mathbb{R}^{n}$ Therefore,

$$
\|T f\|_{2}=\left\|(T f)^{\wedge}\right\|_{2} \leq\|\hat{k}\|_{\infty}\|\hat{f}\|_{2}=\|\hat{k}\|_{\infty}\|f\|_{2}
$$

Step 2. Fix a radial function $\phi \epsilon c_{o}^{\infty}$ with $\phi(x)=1$ for $|x| \leq \frac{1}{2}$ and $\phi(x)=0$ for $|x| \geq 1$. For $\epsilon>0$, we define $k_{\epsilon}(x)=k(x)(1-\phi(x / \epsilon))$ and $T_{\epsilon} f=k_{\epsilon} * f$. Then we claim that $T_{\epsilon}$ is bounded on $L^{2}$ uniformly in $\epsilon$. To see this, we observe that

$$
\begin{aligned}
\left(T_{\epsilon} f\right) & =\hat{f} \hat{k}_{\epsilon}=\hat{f}(k-k \phi(x / \epsilon)) \\
& =\hat{f} \hat{k}-\hat{f}(\hat{k} *(\phi(x / \epsilon)))
\end{aligned}
$$

which gives

$$
\begin{aligned}
\left\|T_{\epsilon} f\right\|_{2} & =\left\|\left(T_{\epsilon} f\right)\right\|_{2} \leq\|\hat{f}\|_{2}\|\hat{k}\|_{\infty}\left\{1+\epsilon^{n} \int|\hat{\phi}(\epsilon \xi)| d \xi\right\} \\
& =\|f\|_{2}\|\hat{k}\|_{\infty}\left\{1+\|\hat{\phi}\|_{1}\right\}
\end{aligned}
$$

139 Step 3. $T_{\epsilon}$ is of weak type $(1,1)$ uniformly in $\epsilon$. The proof of this is more involved and will be given later.

Step 4. By steps 2and 3, using Marcinkiewicz, we get that $T_{\epsilon}$ is bounded on $L^{p}$ for $1<p<2$ uniformly in $\epsilon$.

Step 5. $T_{\epsilon}$ is bounded on $L^{p}$, for $2<p<\infty$, uniformly in $\epsilon$. Indeed, for $f, g \in C_{o}^{\infty}$,

$$
\int\left(T_{\epsilon} f\right)(x) g(x) d x=\int f(x)\left(\tilde{T}_{\epsilon} g\right)(x) d x
$$

with $\tilde{T}_{\epsilon} g=\tilde{k}_{\epsilon} * g, \tilde{k}_{\epsilon}=k_{\epsilon}(-x)=k(-x)(1-\phi(x / \epsilon))$. Since $\tilde{k}_{\epsilon}$ satisfies the same conditions as $k_{\epsilon}$, we see that $\tilde{T}_{\epsilon}$ is bounded on $L^{q}$ for $1<q<2$ uniformly in $\epsilon$. So if $(1 / p)+(1 / q)=1,1<q<2$,

$$
\left\|T_{\epsilon} f\right\|_{p}=\sup _{g \in C_{o}^{\infty}} \frac{\left|\int\left(T_{\epsilon} f\right) g\right|}{\|g\|_{q}} \leq \sup _{g \in C_{o}^{\infty}} \frac{\left\|\tilde{T}_{\epsilon} g\right\|_{q}}{\|g\|_{q}}\|f\|_{p}=c\|f\|_{p}
$$

Step 6. If $f \epsilon C_{o}^{\infty}, T_{\epsilon} f$ converges to $T f$ in the $L^{p}$ norm, $1 \leq p \leq \infty$ as $\epsilon$ tends to 0 . Since $\phi$ is radial and $\mu_{k}=0$,

$$
\int_{|x|=r} \phi(x / \epsilon) k(x) d \sigma(x)=c \int_{|x|=r} k(x) d \sigma(x)=0
$$

so, for $\epsilon \leq 1$, we have
$\left(T_{\epsilon} f\right)(x)=\int_{|y| \leq 1}(f(x-y)-f(x)) k(y)(1-\phi(y / \epsilon)) d y+\int_{|y|>1} f(x-y) k(y) d y$
and hence $\quad\left(T_{\epsilon} f\right)(x)-(T f)(x)=\int_{|y| \leq \epsilon} f(x-y)-f(x) k(y) \phi(y / \epsilon) d y$.

Now $\operatorname{supp}\left(T_{\epsilon} f-T f\right) \subset\{x: d(x, \operatorname{supp} f) \leq \epsilon\} \subset A$, a fixed compact set. Since, on compact sets, the uniform norm dominates all $L^{p}$ norms, it suffices to show that $\left(T_{\epsilon} f-T f\right)$ converges to 0 uniformly on $A$. But

$$
\begin{aligned}
\left\|T_{\epsilon} f-f\right\|_{\infty} \leq\|\operatorname{grad} f\|_{\infty} & \int_{|y| \leq \epsilon} c\left|y\left\|\left.y\right|^{-n} d y \leq c^{\prime}\right\| \operatorname{grad} f \|_{\infty}\right. \\
& \int_{0}^{\epsilon} r^{1-n} r^{n-1} d r=c^{\prime \prime} \epsilon \rightarrow 0, \text { as } \epsilon \rightarrow 0
\end{aligned}
$$

Step 7. If $f \epsilon L^{p}$ and $\eta>0$, choose $g \epsilon C_{o}^{\infty}$ with $\|g-f\|_{p}<\eta$.
Then $\left\|T_{\epsilon} f-T_{\delta} f\right\|_{p} \leq\left\|T_{\epsilon}(f-g)\right\|_{p}+\left\|T_{\epsilon} g-T_{\delta} g\right\|_{p}+\left\|T_{\delta}(g-f)\right\|_{p} \leq$ $2 c \eta+\left\|T_{\epsilon} g-T_{\delta} g\right\|_{p}$.

Since $\eta$ is arbitrary and $\left\|T_{\epsilon} g-T_{\delta} g\right\|_{p}$ converges to 0 as $\epsilon, \delta$ tend to 0 by step 6, we see that $\left(T_{\epsilon} f\right)$ is Cauchy in the $L^{p}$ norm. Setting

$$
T f=\lim _{\epsilon \rightarrow 0} T_{\epsilon} f,\|T f\|_{p} \leq c\|f\|_{p}
$$

Thus the theorem is proved modulo Step 3
Let us now proceed to the proof of Step 3. First, we need a lemma which will be used in the proof.

Lemma 5.17. Suppose $F \epsilon L^{1}, F \geq 0$ and $\alpha>0$. Then there exists a sequence $\left(Q_{k}\right)$ of closed cubes with sides parallel to the coordinate axes and disjoint interiors such that
a) $\alpha<\frac{1}{\left|Q_{k}\right|} \int_{Q_{k}} F \leq 2^{n} \alpha$ for all $k$,
b) If $\Omega=\bigcup_{1}^{\infty} Q_{k}$, then $|\Omega| \leq(1 / \alpha)\|F\|_{1}$,
c) $F(x) \leq \alpha$ for a.e. $x \notin \Omega$.

Proof. Let $r=\left(\|F\|_{1} / \alpha\right)^{1 / n}$ and for $j=1,2, \ldots$ let $\mathscr{Q}_{j}$ be the collection of closed cubes of side length $r / 2^{j}$ and vertices in $\left(r / 2^{j}\right) \mathbb{Z}^{n}$. Our sequence will be constructed in the following way. Put those cubes $Q \in \mathscr{Q}_{1}$ in the sequence which satisfy $\alpha<\frac{1}{|Q|} \int_{Q} F$. Then

$$
\frac{1}{|Q|} \int_{Q} F \leq \frac{1}{|Q|}\|F\|_{1}=\left(2^{n} \alpha /\|F\|_{1}\right)\|F\|_{1}=2^{n} \alpha
$$

so that the first condition is satisfied.
Put those cubes $Q \epsilon \mathscr{Q}_{2}$ into the sequence which are not contained in one of the previously accepted cubes and satisfy

$$
\alpha<\frac{1}{|Q|} \int_{Q} F
$$

Inductively, put those $Q \epsilon \mathscr{Q}_{j}$ which are not contained in one of the previously accepted cubes and satisfy $\alpha<\frac{1}{|Q|} \int_{Q} F$. If $Q \epsilon \mathscr{Q}_{j}$ is in the sequence and $Q^{\prime}$ is the cube in $\mathscr{Q}_{j-1}$ containing $Q$ then

$$
\frac{1}{|Q|} \int_{Q} F \leq \frac{1}{|Q|} \int_{Q}, F=\frac{2^{n}}{\left|Q^{\prime}\right|} \int_{Q^{\prime}} F
$$

Since $Q \subset Q^{\prime}, Q^{\prime}$ cannot be in the sequence and hence

$$
\frac{1}{\left|Q^{\prime}\right|} \int_{Q^{\prime}} F \leq \alpha \text { which then yields } \frac{1}{|Q|} \int_{Q} F \leq 2^{n} \alpha
$$

Thus the condition $(a)$ is satisfied. Also,

$$
|\Omega|=\sum_{1}^{\infty}\left|Q_{k}\right| \leq \frac{1}{\alpha} \int_{Q_{k}} F \leq(1 / \alpha)\|F\|_{1} ; \text { so }(b) \text { follows }
$$

Finally, by the Lebesgue differentiation theorem,

$$
\lim _{x \in Q \epsilon \mathscr{Q}_{j} \rightarrow \infty} \frac{1}{|Q|} \int_{Q} F=F(x) \text { for a.e. } x .
$$

So, if $x \notin \Omega, \frac{1}{|Q|} \int_{Q} F \leq \alpha$ for all those $Q^{\prime}$ s and hence $F(x) \leq \alpha$ a.e. on $\mathbb{R}^{n} \backslash \Omega$.

Coming back to the proof of Step3 given $f \epsilon L^{1}$ and $\alpha>0$, let $\left(Q_{k}\right)$ be the sequence of cubes as in the lemma with $F=|f|$. We write

$$
\begin{aligned}
f & =g+\sum_{k=1}^{\infty} b_{k} \text { with } \\
b_{k}(x) & =\left\{\begin{array}{l}
f(x)-\frac{1}{Q_{k} \mid} \int_{Q_{k}} f(y) d y, \text { for } x \in Q_{k} \\
0, \text { otherwise }
\end{array}\right. \\
\text { and } g(x) & =\left\{\begin{array}{l}
\frac{1}{\left|Q_{k}\right|} \int_{Q_{k}} f(y) d y, \text { for } x \in Q_{k} \\
f(x), \text { for } x \notin \Omega .
\end{array}\right.
\end{aligned}
$$

Now me make the following observations :
a) Supp $b_{k} \subset Q_{k}, \int b_{k}=0$ and

$$
\begin{equation*}
\sum_{1}^{\infty}\left\|b_{k}\right\|_{1} \leq \sum_{1}^{\infty} 2 \int_{Q_{k}}|f| \leq 2 \sum_{1}^{\infty} 2^{n} \alpha\left|Q_{k}\right| \leq 2^{n+1}\|f\|_{1} \tag{5.18}
\end{equation*}
$$

b) $|g(x)| \leq 2^{n} \alpha$ for $x \in \Omega$ and $|g(x)| \leq \alpha$ for a.e. $x \notin \Omega$.

Therefore, $|g(x)| \leq 2^{n} \alpha$ a.e. and

$$
\begin{align*}
\|g\|_{2}^{2} & =\int_{\Omega}|g|^{2}+\int_{\mathbb{R}^{n}-\Omega}|g|^{2}  \tag{5.19}\\
& \leq\left(2^{n} \alpha\right)^{2}|\Omega|+\alpha \int_{\mathbb{R}^{n}-\Omega}\|f\| \\
& \leq\left(2^{2 n}+1\right) \alpha\|f\|_{1} .
\end{align*}
$$

We put $\sum_{1}^{\infty} b_{k}=b$ so that $T_{\epsilon} f=T_{\epsilon} g+T_{\epsilon} b$ and

$$
\left|\left\{\left|T_{\epsilon} f\right|>\alpha\right\}\right| \leq\left|\left\{\left|T_{\epsilon} g\right|>\alpha / 2\right\}\right|+\left|\left\{\left|T_{\epsilon} b\right|>\alpha / 2\right\}\right| .
$$

We shall show that both terms on the right are dominated by $\|f\|_{1 / \alpha}$ uniformly in $\epsilon$.

To estimate the first term on the right, we use Chebyshev inequality, Step 2and the estimate (5.19) obtaining

$$
\left.\left|\left\{\left|T_{\epsilon} g\right|>\alpha / 2\right\}\right| \leq(2 / \alpha)\left\|T_{\epsilon} g\right\|_{2}\right)^{2} \leq c\left(\|g\|_{2} / \alpha\right)^{2} \leq c_{1}\left(\|f\|_{1} \alpha\right) .
$$

To estimate the second term, let $y_{k}$ be the center of the cube $Q_{k}$ and $\tilde{Q}_{k}$ be the cube centred at $y_{k}$ but with length side $2 \sqrt{n}$ times that of $Q_{k}$. We put $\bigcup_{1}^{\infty} \tilde{Q}_{k}=\tilde{\Omega}$. Then

$$
\begin{aligned}
|\tilde{\Omega}| \leq \sum_{1}^{\infty}\left|\tilde{Q}_{k}\right|=(2 \sqrt{n})^{n} \sum_{1}^{\infty} & \left|Q_{k}\right| \leq(2 \sqrt{n})^{n}\|f\|_{1} / \alpha=c\left(\|f\|_{1} / \alpha\right) . \\
\mid\left\{\left|T_{\epsilon} b\right|>\alpha / 2\right\} & \leq|\tilde{\Omega}|+\left|\left\{\left|T_{\epsilon} b\right|>\alpha / 2\right\} \backslash \tilde{\Omega}\right| \\
& \leq c\left(\|f\|_{1} / \alpha\right)+\left|\left\{\left|T_{\epsilon} b\right|>\alpha / 2\right\} \backslash \tilde{\Omega}\right|,
\end{aligned}
$$

So

144 and it suffices to estimate $\mid\left\{\left|T_{\epsilon}\right|>\alpha / 2\right\} / \tilde{\Omega}$. Since $\int_{Q_{k}} b=0$,

$$
T_{\epsilon} b(x)=\int k_{\epsilon}(x-y) b(y) d y
$$

$$
\begin{aligned}
& =\sum_{1}^{\infty} \int_{Q_{k}} k_{\epsilon}(x-y) b(y) d y \\
& =\sum_{1}^{\infty} \int_{Q_{k}}\left(k_{\epsilon}(x-y)-k_{\epsilon}\left(x-y_{k}\right)\right) b(y) d y
\end{aligned}
$$

Therefore, we have

$$
\begin{aligned}
\mid\left\{\left|T_{\epsilon} b\right|\right. & >\alpha / 2\} \backslash \tilde{\Omega} \mid \\
& \leq(2 / \alpha) \int_{\mathbb{R}^{n} \backslash \tilde{\Omega}}\left|T_{\epsilon} b\right| d x \\
& \leq(2 / \alpha) \sum_{1}^{\infty} \int_{\mathbb{R}^{n} \backslash \tilde{\Omega}} \int_{Q_{k}}\left|\left(k_{\epsilon}(x-y)-k_{\epsilon}\left(x-y_{k}\right)\right)\right||b(y)| d y d x \\
& \leq(2 / \alpha) \sum_{1}^{\infty} \int_{Q_{k}} \int_{\mathbb{R}^{n} \backslash \tilde{Q}_{k}}\left|\left(k_{\epsilon}(x-y)-k_{\epsilon}\left(x-y_{k}\right)\right)\right||b(y)| d x d y
\end{aligned}
$$

We now claim that
$\int_{\mathbb{R}^{n} \backslash \tilde{Q}_{k}}\left|\left(k_{\epsilon}(x-y)-k_{\epsilon}\left(x-y_{k}\right)\right)\right| d x \leq c$ independent of $k$ and $\epsilon$ for $y \epsilon Q_{k}$.
Accepting this claim, by the estimate (5.18, we have

$$
\left|\left\{\left|T_{\epsilon}\right|>\alpha / 2\right\} \backslash \tilde{\Omega}\right| \leq(2 / \alpha) c \sum_{1}^{\infty} \int_{Q_{k}}|b(y)| d y \leq 2^{n+2} c\left(\|f\|_{1} / \alpha\right)
$$

Thus $\left|\left\{\left|T_{\epsilon} b\right|>\alpha\right\}\right| \leq c_{o}\left(\|f\|_{1} / \alpha\right)$ which completes the proof of Step (3)

Returning to the claim, we observe that if $x \in \mathbb{R}^{n} \backslash \tilde{Q}_{k}$ and $y \in Q_{k}$, then $\left|x-y_{k}\right| \geq 2\left|y-y_{k}\right|$.

So, if we set $z=x-y_{k}, w=y-y_{k}$ we must show that

$$
\int_{|z|>2|w|}\left|\left(k_{\epsilon}(z-w)-k_{\epsilon}(z)\right)\right| d z \leq c
$$

independent of $w$ and $\epsilon$.


First consider the case $\epsilon=1$. Now $k_{1}(z)$ is a $C^{\infty}$ function which is homogeneous of degree $-n$ for large $z$. Therefore, $\left|\operatorname{gradk}_{1}(z)\right| \leq$ $c^{\prime \prime}|z|^{-n-1}$. By the mean value theorem,

$$
\begin{aligned}
\left|k_{1}(z-w)-k_{1}(z)\right| & \leq c^{\prime}|w| \sup _{0<t<1}|z-t w|^{-n-1} \\
& \leq c^{\prime \prime}|w||z|^{-n-1} \text { for }|z|>2|w|
\end{aligned}
$$

So

$$
\begin{aligned}
& \int_{|z|>2|w|}\left|\left(k_{1}(z-w)-k_{1}(z)\right)\right| d z \\
& \quad \leq c^{\prime \prime} \int_{|z|>2|w|}|w \| z|^{-n-1} d z \\
& \quad \leq c^{\prime \prime \prime}|w| \int_{2|w|}^{\infty} r^{-2} d r=c .
\end{aligned}
$$

Now for general $\epsilon$,

$$
k_{\epsilon}(z)=(1-\phi(z / \epsilon)) k(z / \epsilon) \epsilon^{-n} \text { by the homogeneity of } k .
$$

Therefore, if we set $z^{\prime}=\epsilon^{-1} z$ and $w^{\prime}=\epsilon^{-1} w$, we see that

$$
\int_{|z|>2|w|}\left|\left(k_{\epsilon}(z-w)-k_{\epsilon}(z)\right)\right| d z=\int_{\left|z^{\prime}\right|>2\left|w^{\prime}\right|}\left|\left(k_{1}\left(z^{\prime}-w^{\prime}\right)-k_{1}\left(z^{\prime}\right)\right)\right| d z^{\prime}
$$

which is bounded by a constant, by the result for $\epsilon=1$. Hence the claim above is proved.

To complete the picture, we should observe that operators of type

0 are not bounded on $L^{1}$ (and hence, by duality, not bounded on $L^{\infty}$ ). Indeed, if $T$ is an operator of type 0 with kernel $k,(T f)^{\wedge}=\hat{k} \hat{f}$. Since $\hat{f}$ is homogeneous of degree 0 , it has a discontinuity at 0 (unless $k=c \delta$ ). Thus if $f \epsilon L^{1},(T f)^{\wedge}$ is not continuous at 0 whenever $\hat{f}(0) \neq 0$ and this implies that $T f$ is not in $L^{1}$.

This reflects the fact that if $\int f=\hat{f}(0) \neq 0$, then $T f$ will not be integrable near $\infty$, because $k$ is itself not integrable at $\infty$. However, there are also problems with the local integrability of $T f$ caused by the singularity of $k$ at the origin. In fact, let $\phi \epsilon C_{o}^{\infty}$ be a radial function such that $\phi=1$ near 0 , and set $S f=f *(\phi k)$. Then the argument used to prove Theorem 5.16 shows that $S$ is bounded on $L^{p}$ for $1<p<\infty$. In this case $(\phi k)^{\wedge}=\hat{\phi} * \hat{k}$ is in $C^{\infty}$, but still $S$ is not bounded on $L^{1}$.

This follows from the following general fact.
Proposition 5.20. If $k \in S^{\prime}$ and the operator $f \rightarrow k * f$ is bounded on $L^{1}$, then $k$ is necessarily a finite Borel measure.

Proof. Choose $\phi \epsilon C_{o}^{\infty}$ with $\int \phi=1$ and put $\phi_{\epsilon}(x)=\epsilon^{-n} \phi(x / \epsilon)$. Then $\left\|\phi_{\epsilon}\right\|_{1}$ is independent of $\epsilon$, so $\left\|\phi_{\epsilon} * k\right\|_{1} \leq c$. Therefore there exists a sequence $\epsilon_{k}$ tending to 0 such that $\phi_{\epsilon_{k}} * k$ converges to a finite Borel measure $\mu$ in the weak* topology of measures and hence $\phi_{\epsilon_{k}} * k$ converges to $\mu$ in $S^{\prime}$ also. On the other hand, since $\left(\phi_{\epsilon_{k}}\right)$ is an approximate identity, $\phi_{\epsilon_{k}} * k$ converges to $k$ in $S^{\prime}$.

Hence $\mu=k$ and the proposition is proved.
It is easy to see that kernels of type 0 are not measures even when truncated away from the origin, as their total variation in any neighbourhood of 0 is infinite. One can also see directly that they do not define bounded functionals on $C_{o}$.

Exercise. Let $k(x)=1 / x$ on $\mathbb{R}$ and let $f$ be a continuous compactly supported function such that $f(x)=(\log x)^{-1}$ for $0<x<\frac{1}{2}$ and $f(x)=$ 0 for $x \leq 0$. Show that

$$
\lim _{x \rightarrow 0-}(f * P V(k))(x)=\infty .
$$

Generalise this to kernels of type 0 on $\mathbb{R}^{n}, n>1$.
This example also shows that operators of type 0 do not map continuous functions into continuous functions. However, they do preserve Lipschitz or Hölder continuity, as we shall now see.
Definition 5.21. For $0<\alpha<1$, we define $|f|_{\alpha}=\sup _{x, y} \frac{|f(x+y)-f(x)|}{|y|^{\alpha}}$

$$
\Lambda_{\alpha}=\left\{f:\|f\|_{\Lambda_{\alpha}}={ }^{d e f}\|f\|_{\infty}+|f|_{\alpha}<\infty\right\} .
$$

$\Lambda_{\alpha}$ is called Lipschitz class of order $\alpha$.
Remark 5.22. The definition makes perfectly good sense for $\alpha=1$ (When $\alpha>1$ it is an easy exercise to show that if $|f|_{\alpha}<\infty$ then $f$ is constant ). However, we shall not use this definition for $\alpha=1$, because the theorems we wish to prove are false in this case.

We are going to prove, essentially, that operators of type 0 are bounded on $\Lambda_{\alpha}$. However, if $k$ is a kernel of type 0 and $f \epsilon \Lambda_{\alpha}$, the integral defining $k * f$ will usually diverge $f$ need not decay at $\infty$. Consequently, we shall work instead with $\Lambda^{\alpha} \cap L^{p}(1<p<\infty)$, concerning which we have the following useful result.
Proposition 5.23. If $f \in L^{p}, 1 \leq p<\infty$ and $|f|_{\alpha}<\infty$, then $f \in \Lambda_{\alpha}$ and $\|f\|_{\infty} \leq c\left(\|f\|_{p}+|f|_{\alpha}\right)$. Consequently, $L^{p} \cap \Lambda_{\alpha}$ is a Banach space with norm $\|f\|_{p}+|f|_{\alpha}$.

Proof. Let

$$
A_{x}=\left(\frac{|f(x)|}{2|f|_{\alpha}}\right)^{1 / \alpha} \text { for } x \in \mathbb{R}^{n} .
$$

Then $|f(y)| \geq|f(x)| / 2$ for all $y$ such that $|x-y| \leq A_{x}$. So

$$
\begin{aligned}
\int|f(y)|^{p} d y & \geq \int_{|x-y| \leq A_{x}}|f(y)|^{p} d y \geq|f(x)|^{p} 2^{-p} c^{\prime} A_{x}^{n} \\
& =c^{\prime \prime}|f(x)|^{p+(n / \alpha)}|f|_{\alpha}^{-n / \alpha}
\end{aligned}
$$

or

$$
|f(x)|^{1+(n / \alpha p)} \leq c^{\prime \prime \prime}\|f\|_{p}|f|_{\alpha}^{n / a l p h a p} .
$$

Since this is true for all $x$, setting $\theta=n / p \alpha$ we have

$$
\begin{aligned}
\|f\|_{\infty} & \leq c^{\prime \prime \prime}\|f\|_{p}^{1 /(1+\theta)}|f|_{\alpha}^{\theta /(1+\theta)} \\
& \leq c\left(\|f\|_{p}+|f|_{\alpha}\right)
\end{aligned}
$$

Theorem 5.24. Operators of type 0 are bounded on $\Lambda_{\alpha} \cap L^{p}(0<\alpha<$ $1,<1<p<\infty)$.

Proof. Let $T: f \rightarrow K * f$ be an operator of type 0 . Since we know that $\|T f\|_{p} \leq c_{p}\|f\|_{p}$, by proposition 5.23, it will suffice to show $|T f|_{\alpha} \leq$ $c_{\alpha}|f|_{\alpha}$ for $0<\alpha<1, f \epsilon L^{p} \cap \Lambda_{\text {alpha. }}$. As in the proof of Theorem 5.16 we may assume that $K=P V(k)$ and identify $K$ with $k$.

Given $y \in \mathbb{R}^{n} \backslash\{0\}$ and $f \epsilon L^{p} \cap \Lambda_{\alpha}$ define

$$
\begin{aligned}
& g(x)=\int_{|z| \leq 3|y|} k(z) f(x-z) d z \\
& h(x)=\int_{|z| \leq 3|y|} k(z) f(x-z) d z
\end{aligned}
$$

so that $T f=g+h$. Since $\mu_{k}=0$, we have

$$
\begin{aligned}
|g(x)| & =\left|\int_{|z| \leq 3|y|} k(z)(f(x-z)-f(x)) d z\right| \\
& \leq \int_{|z| \leq 3|y|} c|z|^{-n}|f|_{\alpha}|z|^{\alpha} d z \leq c_{1}|f|_{\alpha}|y|^{\alpha}
\end{aligned}
$$

Since this is true for all $x$,

$$
|g(x+y)-g(x)| \leq 2 c_{1}|f|_{\alpha}|y|^{\alpha} .
$$

Next

$$
h(x+y)=\lim _{\eta \rightarrow \infty} \int_{3|y|<|z|<\eta} k(z)(f(x+y-z)-f(x)) d z
$$

$$
=\lim _{\eta \rightarrow \infty} \int_{3|y|<|z+y|<\eta} k(z+y)(f(x-z)-f(x)) d z
$$

Therefore,
$h(x+y)-h(x)=\lim _{\eta \rightarrow \infty} \int_{3|y|<|z|<\eta}(k(z+y)-k(z))(f(x-z)-f(x)) d z+\epsilon_{1}+\epsilon_{2}$
where $\epsilon_{1}$ and $\epsilon_{2}$ are errors coming from difference between the regions of integration.
$\epsilon_{1}$ is the error coming from the difference between the regions $|z|<\eta$ and $|z+y|<\eta$.

The symmetric difference between these two regions is contained in the annulus $\eta-|y|<|z|<\eta+|y|$.


If $z$ is in this region and $\eta \gg|y|$ then $|z| \approx|z+y| \approx \eta$ so that

$$
\begin{array}{r}
\left|\epsilon_{1}\right| \leq c^{\prime} \int_{\eta-|y|<|z|<\eta+|y|} 2 \eta^{-n}\|f\|_{\infty} d z \leq c^{\prime \prime} \eta^{-n}\|f\|_{\infty}\left((\eta+|y|)^{n}-(\eta-|y|)^{n}\right) \\
=0\left(\eta^{-1}\right) \rightarrow 0 \text { as } \eta \rightarrow \infty
\end{array}
$$

The term $\epsilon_{2}$ comes from the symmetric difference of the regions $|z|>$ $3|y|$ and $|z+y|>3|y|$ which is contained in the annulus $2|y|<|z|<4|y|$.

In this region $|z+y| \approx|z| \approx|y|$.
Therefore

$$
\begin{aligned}
\left|\epsilon_{2}\right| & \leq c^{\prime \prime \prime} \int_{2|y|<|z|<4|y|}|y|^{-n+\alpha}|f|_{\alpha} d z= \\
& =c_{2}^{\prime}|f|_{\alpha}|y|^{-n+\alpha}\left((4|y|)^{n}-(2|y|)^{n}\right) \\
& =c_{2}|f|_{\alpha}|y|^{\alpha} .
\end{aligned}
$$

Finally coming to the main term, we have

$$
\begin{aligned}
|k(z+y)-k(z)| & \leq|y| \sup _{0<t<1}|\operatorname{grad} k(z+t y)| \\
& \leq|y| \sup _{0<t<1}|z+t y|^{-n-1} \\
& \leq c_{o}|y||z|^{-n-1} \text { for }|z| \geq 3|y| .
\end{aligned}
$$

Hence, since $\alpha<1$ so that $-n-1+\alpha<-n$,

$$
\begin{aligned}
& \left|\int_{3|y|<|z|<\eta}(k(z+y)-k(z))(f(x-z)-f(x)) d z\right| \\
& \quad \leq c_{o} \int_{3|y|<|z|<\eta}|y||z|^{-n-1}|f|_{\alpha}|z|^{\alpha} d z \\
& \quad \leq c_{o} \int_{3|y|<|z|}|y||z|^{-n-1+\alpha}|f|_{\alpha} d z \\
& \quad \leq c_{3}|y||f|_{\alpha}|y|^{\alpha-1}=c_{3}|y|^{\alpha}|f|_{\alpha}
\end{aligned}
$$

Therefore,

$$
\frac{|T f(x+y)-T f(x)|}{|y|^{\alpha}} \leq c|f|_{\alpha} \text { and consequently }|T f|_{\alpha} \leq c|f|_{\alpha}
$$

For kernels of positive type, we have the following result.
Theorem 5.25. Suppose $0<\lambda<n, 1<p<n / \lambda<q<Q$ where $Q=\infty$ if $\lambda \leq 1, Q=n /(\lambda-1)$, if $\lambda<1$. Let $1 / r=(1 / p)-(\lambda / n)$ and $\alpha=\lambda-(n / q) .($ Thus $r<\infty$ and $0<\alpha<1)$. Thus operators of type $\lambda$ are bounded from $L^{p} \cap L^{q}$ into $L^{r} \cap \Lambda_{\alpha}$.

Proof. Let $T f=k * f$ be an operator of type $\lambda$. By Corollary $5.15 T$ is bounded from $L^{p}$ to $L^{r}$, so by proposition 5.23 it is enough to show that $|T f|_{\alpha} \leq c\|f\|_{q}$.

We have $(T f)(x)=\int f(x-z) k(z) d z$

$$
\begin{aligned}
(T f)(x+y) & =\int f(x+y-z) k(z) d z \\
& =\int f(x-z) k(z+y) d z
\end{aligned}
$$

so that

$$
\begin{aligned}
(T f)(x+y)-T f(x)= & \int f(x-z)(k(z+y)-k(z)) d z \\
= & \int_{|z| \leq 2|y|} f(x-z)(k(z+y)-k(z)) d z \\
& +\int_{|z|>2|y|} f(x-z)(k(z+y)-k(z)) d z .
\end{aligned}
$$

If $q^{\prime}$ is the conjugate exponent of $q$,

$$
\begin{aligned}
& \left|\int_{|z| \leq 2|y|} f(x-z)(k(z+y)-k(z)) d z\right| \\
& \quad \leq\|f\|_{q}\left\{\left(\int_{|w| \leq 3|y|} \mid k(w)^{q^{\prime}} d w\right)^{1 / q^{\prime}}+\left(\int_{|z| \leq 2|y|} \mid k(z)^{q^{\prime}} d z\right)^{1 / q^{\prime}}\right) \\
& \quad \leq 2| | f \|_{q}\left(\int_{|z| \leq 3|y|}|k(z)|^{q^{\prime}} d z\right)^{1 / q^{\prime}} \\
& \quad \leq c_{1}\|f\|_{q}\left(\int_{|z| \leq 3|y|}|z|^{(\lambda-n) q^{\prime}} d z\right)^{1 / q^{\prime}} \\
& \quad \leq c_{1}^{\prime}\|f\|_{q|y|^{-n+\lambda+\left(n / q^{\prime}\right)}=c_{1}^{\prime}\|f\|_{q}|y|^{\alpha}}
\end{aligned}
$$

by the definition of $\alpha$
For the second integral, we estimate $k(z+y)-k(z)$ by the mean value theorem:

$$
|k(z+y)-k(z)| \leq|y| \sup _{0<t<1}|\operatorname{grad} k(z+t y)|
$$

$$
\leq c|y||z|^{\lambda-n-1} \text { for }|z| \geq 2|y|
$$

Thus

$$
\begin{aligned}
& \left|\int_{|z|>2|y|} f(x-z)(k(z+y)-k(z)) d z\right| \\
& \quad \leq c\|f\|_{q}\left\{\int_{|z|>2|y|}\left(|y \| z|^{\lambda-n-1}\right)^{q^{\prime}} d z\right\}^{1 / q^{\prime}} \\
& \quad \leq c_{2}^{\prime}\|f\|_{q}|y \| y|^{\lambda-n-1+\left(n / q^{\prime}\right)} \operatorname{since}(\lambda-n-1) q^{\prime}<-n \\
& \quad=c_{2}^{\prime}\|f\|_{q}|y|^{\alpha} .
\end{aligned}
$$

Hence

$$
\frac{|T f(x+y)-T f(x)|}{|y|^{\alpha}} \leq c\|f\|_{q} \text { which gives }|T f|_{\alpha} \leq c\|f\|_{q} .
$$

Remark 5.26. As in the preceding theorem, the reason for taking the domain of $T$ to be $L^{p} \cap L^{q}$ instead of just $L^{q}$ is that the integral defining $T f$ will usually diverge when $f$ is merely in $L^{q}$.

Nonetheless, the point of these theorems is that operators of type 0 are in essence, bounded from $L^{q}$ to $\Lambda_{\alpha}$ (for appropriate $q$ and $\alpha$ ). To make this precise without losing simplicity, one can observe that operators of type 0 map $\Lambda_{\alpha} \cap E^{\prime}$ into $\Lambda_{\alpha}$ while operators of type $\lambda>0$ map $L^{q} \cap E^{\prime}$ into $\Lambda_{\alpha}$.

We now introduce spaces of functions whose derivatives upto a certain order are in $L^{p}$ or $\Lambda_{\alpha}$.

Definition 5.27. Suppose $1 \leq p \leq \infty$ and $k$ is a positive integer. We define

$$
L_{k}^{p}=\left\{f: D^{\beta} f \epsilon L^{p} \text { for } 0 \leq|\beta| \leq k\right\} .
$$

We equip $L_{k}^{p}$ with the norm $\|f\|_{k, p}=\sum_{|\beta| \leq k}\left\|D^{\beta} f\right\|_{p}$. (Thus $L_{k}^{2}=H_{k}$ in the notation of Chapter 3.

Definition 5.28. Suppose $k$ is a positive integer and $k<\alpha<k+1$. We define

$$
\Lambda_{\alpha}=\left\{f: D^{\beta} f \epsilon \Lambda_{\alpha-k} f o r 0 \leq|\beta| \leq k\right\}
$$

We equip $\Lambda_{\alpha}$ with the norm $\|f\|_{\Lambda_{\alpha}}=\sum_{|\beta| \leq k}\left\|D^{\beta} f\right\|_{\Lambda_{\alpha-k}}$.

## Remarks 5.29.

154 (i) $f \epsilon \Lambda_{\alpha}$ if and only if $D^{\beta} f$ is bounded and continuous for $0 \leq|\beta| \leq k$ and $D^{\beta} f \epsilon \Lambda_{\alpha-k}$ for $|\beta|=k$. Indeed, if $|\beta|<k$, for $|y| \leq 1$,

$$
\frac{\left|D^{\beta} f(x+y)-D^{\beta} f(x)\right|}{|y|^{\alpha-k}} \leq \frac{\left|D^{\beta} f(x+y)-D^{\beta} f(x)\right|}{|y|} \leq c \sum_{|v|=|\beta|+1}\left\|D^{v} f\right\|_{\infty}
$$

(by the mean value theorem) and for

$$
|y|>1, \frac{\left|D^{\beta} f(x+y)-D^{\beta} f(x)\right|}{|y|^{\alpha-k}} \leq 2\left\|D^{\beta} f\right\|_{\infty}
$$

This shows that $D^{\beta} f \epsilon \Lambda_{\alpha-k}$ for $0 \leq|\beta| \leq k$ i.e., $f \epsilon \Lambda_{\alpha}$.
(ii) If $k<\alpha<k+1$ then $L_{k}^{p} \cap \Lambda_{\alpha}$ is a Banach space with the norm

$$
\sum_{|\beta| \leq k}\left(\left\|D^{\beta} f\right\|_{p}+\left|D^{\beta} f\right|_{\alpha-k}\right)
$$

This follows from the corresponding fact that $L^{p} \cap \Lambda_{\alpha}$ is a Banach space (Proposition 5.23).

Theorem 5.30. Suppose $0 \leq \lambda<n, 1<p<n / \lambda, 1 / r=(1 / p)-(\lambda / n)$ and $k=0,1,2, \ldots$ Then we have:
a) Operators of type $\lambda$ are bounded from $L_{k}^{p}$ into $L_{k}^{r}$.
b) If $\lambda$ is an integer, $\lambda=0,1,2, \ldots$ and $k<\alpha<k+1$, then operators of type $\lambda$ are bounded from $L_{k}^{p} \cap \Lambda_{\alpha}$ to $L_{k+\lambda}^{r} \cap \Lambda_{\alpha+\lambda}$.

Proof. a) This is an easy consequence of Corollary 5.15 and Theorem 5.16 since convolution commutes with differentiation. In the same way (b) follows Theorem 5.24 when $\lambda=0$.

We now proceed by induction on $\lambda$. Therefore, assume that $\lambda \geq 1$. Suppose $f \epsilon L_{k}^{p} \cap \Lambda_{\alpha}$. Then $\partial_{j}(f * k)=f * \partial_{j} k$ and $\partial_{j} k$ is a kernel of type $\lambda-1$. So, if $1 / s=(1 / p)-(\lambda-1) / n$, by our induction hypothesis $\partial_{j}(f * k) \epsilon L_{k+\lambda-1}^{s} \cap \Lambda_{\alpha+\lambda-1}$. Since $s<r<\infty, L^{s} \cap L^{\infty} \subset L^{r}$, so $\partial_{j}(f *$ $k) \epsilon L_{k+\lambda-1}^{r} \stackrel{(1+\lambda-1}{\cap} \Lambda_{\alpha+\lambda-1}$. Also $f * k \epsilon L^{r} \cap C^{1}$, hence in $\Lambda_{\alpha+\lambda}$ provided it is bounded. But $\partial_{j}(f * k) \in \Lambda_{\alpha+\lambda-1}$ implies that $\partial_{j}(f * k)$ is bounded. By the mean value theorem, we then have

$$
\frac{|(f * k)(x+y)-(f * k)(x)|}{|y|} \leq c .
$$

This together with $f * k \in L^{r}$ implies that $f$ is bounded (by definition 5.21). Hence $f * k \in L_{k+\lambda}^{r} \cap \Lambda_{\alpha+\lambda}$.

The above theorem can be generalised. For example, if $0 \leq \lambda<n$ one can show that operators of type $\lambda$ map $\Lambda_{\alpha} \cap E^{\prime}$ in to $\Lambda_{\alpha+\lambda}$.

Also generalisations of the $L^{p}$ Sobolev spaces $L_{k}^{p}$ can be given for non-integral values of $k$. In fact, a theorem due to CALDERON says that for $1<p<\infty, f \in L_{k}^{p}$ if and only if $\Lambda^{k} f \in L^{p}$. (Here $\Lambda=(1-\Delta)^{1 / 2}$ ).

Therefore, we can define $L_{s}^{p}$ for any real $s$ by
$L_{s}^{p}=\left\{f: \Lambda^{s} f \in L^{p}\right\}$ with the norm $\|f\|_{s, p}=\left\|\Lambda^{s} f\right\|_{p}$.
Then part (b) of the above theorem is still true for $0 \leq \lambda<n, \lambda$ not necessarily an integer in this case.

Refer to E.M. Stein [2].
We will now prove the Sobolev imbedding theorem for $L^{p}$ Sobolev spaces $L_{k}^{p}$ with positive integral $k$. This theorem can also be generalised to $L_{s}^{p}$ for $s$ in $\mathbb{R}$.
Theorem 5.31 (SOBOLEV IMBEDDING THEOREM). Suppose $1<$ $p<\infty$ and $k$ a positive integer. If $k<n / p$, then $L_{k}^{p} \subset L^{r}$ for $1 / r=$ $(1 / p)-(k / n)$ (Hence also $L_{k+j}^{p} \subset L_{j}^{r}$ for any $j$ ). If $k>n / p$, and $\alpha=$ $k-n / p$ is not integer, then $L_{k}^{p} \subset \Lambda_{\alpha}$.
Proof. Let $N$ be the fundamental solution of $\Delta$ given by

$$
N\left(x=\left\{\begin{array}{l}
(2-n)^{-1} \omega_{n}^{-1}|x|^{2-n} \text { for } n \neq 2 \\
(2 \pi)^{-1} \log |x| \text { for } n=2
\end{array}\right.\right.
$$

Then $K_{j}(x)=\partial_{j} N(x)=\omega_{n} x_{j}|x|^{-n}$ ( true for all $n$ ) is a kernel of type 1.

Now, if $f \in L_{k}^{p} \cap E^{\prime}$

$$
f=f * \delta=f * \Delta N=f * \sum_{j=1}^{n} \partial_{j} k_{j}=\sum_{j=1}^{n}\left(\partial_{j} f * k_{j}\right)
$$

Suppose $k=1$. If $1<n / p, \partial_{j} f \epsilon L^{p} \Rightarrow \partial_{j} f * k_{j} \epsilon L^{r}$ where $1 / r=$ $(1 / p)-(1 / n)$, by Theorem 5.25 and hence $f \in L^{r}$. If $1>n / p, \partial_{j} f *$ $k_{j} \epsilon \Lambda_{1-(n / p)}$ by Theorem 5.25 which implies that $f \in \Lambda_{1(n / p)}$. Thus the theorem is true for $k=1$.

For $k>1$, we proceed by induction. Now
$f \epsilon L_{k}^{p} \cap E^{\prime} \Rightarrow f$ and $\partial_{j} f$ are in $L_{k-1}^{p} \cap E^{\prime}$.
Therefore, if $p<n /(k-1)$ and $1 / q=(1 / p)-(k-1) / n$, we have $f$, $\partial_{j} f \epsilon L^{q}$, i.e., $f \in L_{1}^{q}$, while if $p>n /(k-1)$, we have $f, \partial_{j} f \epsilon \Lambda_{k-1-(n / p)}$, i.e., $f \epsilon \Lambda_{k-(n / p)}$. In the second case, we are done, and in the first case, we apply the result for $k=1$ to see that $f$ is in the required space. Finally, it is easy to check that if we keep track of the norm inequalities that are implicit in the above arguments, we obtain

$$
\|f\|_{r} \leq c\|f\|_{k, p} \text { or }\|f\|_{\Lambda_{\alpha}} \leq c\|f\|_{k, p}
$$

as appropriate, for $f \epsilon L_{k}^{p} \cap E^{\prime}$. Since $L_{k}^{p} \cap E^{\prime}$ is clearly dense in $L_{k}^{p}$, the desired result follows immediately.

We can summarise these theorems in an elegant way using the following picture.

For $-n<\alpha<0$ we define $x_{\alpha}=L^{n /|\alpha|}$ and when $\alpha>0, \alpha$ not an integer, we define $x_{\alpha}=\Lambda_{\alpha}$.

(The small circle represent missing spaces)
In this terminology, we have
Theorem 5.32. Operators of type $\lambda$ map $x_{\alpha} \cap E^{\prime}$ into $x_{\alpha+\lambda} O \leq \lambda<n$.
Theorem 5.33 (SOBOLEV IMBEDDING THEOREM). If $D^{\beta} f \epsilon x_{\alpha}$ for $0 \leq|\beta| \leq k$, then $f \in x_{\alpha+k}$.

We now indicate how to fill the gaps in this picture at $\alpha=0,1,2, \ldots$
For $\alpha=1$, we define
$\Lambda_{1}=\{f: f$ is continuous, bounded and

$$
\left.\sup _{x, y} \frac{|f(x+y)+f(x-y)-2 f(x)|}{|y|}<\infty\right\} .
$$

For $k=2,3,4, \ldots$ we define $\Lambda_{k}=\left\{f: D^{\beta} f \epsilon \Lambda_{1}\right.$ for $\left.|\beta| \leq k-1\right\}$. The sudden jump from first differences in the definition of $\Lambda_{\alpha}$ for $\alpha<1$ to the second differences at $\alpha=1$ is less mysterious than it seems at first, because it can be shown that if $0<\alpha<2$ then $f \epsilon \Lambda_{\alpha}$ if and only if $f$ is bounded, continuous and satisfies

$$
\sup _{x, y} \frac{|f(x+y)+f(x-y)-2 f(x)|}{|y|^{\alpha}}<\infty .
$$

To fill the gap at $\alpha=0$ we use the space BMO (" bounded mean oscillation") first introduced by F. JOHN and L. NIRENBERG in 1961, which is defined as follows.

For $f \in L_{\text {loc }}^{1}\left(\mathbb{R}^{n}\right)$, we denote by $m_{E} f$ mean value of $f$ over a measurable set $E \subset \mathbb{R}^{n}$, that is,

$$
m_{E} f=\frac{1}{|E|} \int_{E} f
$$

Let $Q$ denote the collection of all cubes in $\mathbb{R}^{n}$ with sides parallel to the axes.

Definition 5.34. $B M O=\left\{f \in L_{\mathrm{loc}}^{1}\left(\mathbb{R}^{n}\right): \sup _{Q \in \mathscr{Q}} m_{Q}\left(\left|f-m_{Q} f\right|\right)<\infty\right\}$.
Clearly $L^{\infty} \subset B M O$, for, $f \in L^{\infty} \Rightarrow\left|m_{Q} f\right| \leq\|f\|_{\infty}$ for every $Q \epsilon \mathscr{Q}$ and consequently

$$
m_{Q}\left(\left|f-m_{Q} f\right|\right) \leq 2\|f\|_{\infty} .
$$

It can be shown that $B M O \subset L_{\mathrm{loc}}^{q}$ for every $q<\infty$.
If we define $x_{\alpha}=\Lambda_{\alpha}$ for $\alpha=1,2, \ldots$ and $X_{0}=B M O$ then Theorem 5.32]remains valid for all $\alpha_{\epsilon}(-n, \infty)$, except that the Sobolev imbedding theorem for $\alpha=0$ must be slightly modified as follows:

If $D^{\beta} f$ is in the closure of BMO $\cap E^{\prime}$ in BMO for $|\beta| \leq k$ then $f \epsilon \Lambda_{k}$.

WARNING. BMO is not an interpolation space between $L^{p}$ and $\Lambda_{\alpha}$ i.e., it is not true that if $T$ is a linear operator which is bounded on $L^{p}$ for some $p<\infty$ and on $\Lambda_{\alpha}$ some $\alpha>0$, then $T$ is bounded on BMO.

For proofs of the foregoing assertions, see E.M. Stein [2] and also the following papers :
(i) E.M Stein and A. Zygmund: Boundedness of translation invariant operators on Holder spaces and $L^{p}$ spaces, Ann. Math 85 (11967)337-349, and
(ii) C. Fefferman and E.M. Stein : $H^{p}$ spaces of several variables Acta Math. 129, (1972), 137-193.

As we indicated at the beginning of this chapter, the arguments we have developed can be extended in a fairly straightforward way to give estimates for $\psi D O$ with variable coefficients. We conclude by summarising the result in

Theorem 5.35. [Let] $P=p(x, D)$ ba a properly supported $\psi D O$ of order $-\lambda$ on $\Omega$, where $\lambda \geq 0$ and $p \sim \sum_{j=0}^{\infty} p_{j}$ with $p_{j}(x, \xi)$ homogeneous of degree $\lambda$ - j for large $|\xi|$. Then P maps $L_{k}^{p}(\Omega$, loc $)$ into $L_{k+\lambda}^{p}(\Omega$, loc $)$ for $1<p<\infty$, and in the terminology of Theorem $5.32 P$ maps $x_{\alpha}(\Omega$, loc $)$ into $x_{\alpha+\lambda}(\Omega$, loc $)$ for $-n<\alpha<\infty$.

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