# CHARACTERIZATION OF EIGENFUNCTIONS BY BOUNDEDNESS CONDITIONS 

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#### Abstract

Suppose $\left\{f_{k}(x)\right\}_{k=-\infty}^{\infty}$ is a sequence of functions on $\mathbb{R}^{n}$ with $\Delta f_{k}=f_{k+1}$ (where $\Delta$ is the Laplacian) that satisfies the growth condition: $\left|f_{k}(x)\right| \leq M_{k}(1+|x|)^{a}$ where $a \geq 0$ and the constants have sublinear growth $\frac{M_{k}}{k} \rightarrow 0$ as $k \rightarrow \pm \infty$. Then $\Delta f_{0}=-f_{0}$. This characterizes eigenfunctions $f$ of $\Delta$ with polynomial growth in terms of the size of the powers $\Delta^{k} f,-\infty<k<\infty$. It also generalizes results of Roe (where $a=0, M_{k}=M$, and $n=1$ ) and Strichartz (where $a=0, M_{k}=M$ for $n$ ). The analogue holds for formally self-adjoint constant coefficient linear partial differential operators on $\mathbb{R}^{n}$.


1. Introduction. Let $\Delta=\frac{\partial^{2}}{\partial x_{1}^{2}}+\cdots+\frac{\partial^{2}}{\partial x_{n}^{2}}$ be the Laplacian operator on $\mathbb{R}^{n}$. Recently Strichartz [8] has given a characterization of the bounded solutions of $\Delta f=-f$ in terms of bounds on the iterates $\Delta^{k} f$, where $k \in \mathbb{Z}$; more precisely, if $\left\langle f_{k}\right\rangle_{k=-\infty}^{\infty}$ is a doubly infinite sequence of functions on $\mathbb{R}^{n}$ with $\Delta f_{k}=f_{k+1}$ and $\left|f_{k}(x)\right| \leq M$, for some $M$, then $\Delta f_{0}=-f_{0}$. If $n=1$, this is Roe's theorem [6]: If $\left\langle f_{k}\right\rangle_{k=-\infty}^{\infty}$ is a sequence of real-valued functions with $\frac{d}{d x}\left(f_{k}(x)\right)=f_{k+1}(x)$ and $\left|f_{k}(x)\right| \leq M$, then $f_{0}(x)=A \sin (x+\alpha)$. (See also the paper by Burkill [2] and the paper [3] for generalizations of Roe's theorem.)

This does not characterize all solutions of $\Delta f=-f$ on $\mathbb{R}^{n}$ because many are unbounded. For example, let $P_{j}$ be the vector space of all complex-valued polynomials in $x$ and $y$ of degree at most $j$. For reasons of dimension the linear map $p \mapsto p_{x x}+p_{y y}+2 i p_{x}$ from $\mathcal{P}_{j}$ to $\mathcal{P}_{j-1}$ has nontrivial kernel. Let $p$ be a solution of degree $j$ to $p_{x x}+p_{y y}+2 i p_{x}=0$. Then $u(x, y)=p(x, y) e^{i x}$ satisfies $\Delta u=-u$ and has polynomial growth at infinity. In Fourier analysis the functions with polynomial growth are interesting because they are exactly the ones that can be viewed as tempered distributions (i.e., as elements of the dual of the space $\mathcal{S}\left(\mathbb{R}^{n}\right)$ of rapidly decreasing functions (see [4], Chapter 1)). As such they have Fourier transforms that are also tempered distributions. The following gives a characterization of solutions to $\Delta f=-f$ of (at most) polynomial growth in terms bounds on the power $\Delta^{k} f$ for $-\infty<k<\infty$.

Theorem 1. Let $a \geq 0$ and let $\left\langle f_{k}\right\rangle_{-\infty}^{\infty}$ be a sequence of complex-valued functions on $\mathbb{R}^{n}$ that satisfy

$$
\Delta f_{k}=f_{k+1}
$$

and

$$
\left|f_{k}(x)\right| \leq M_{k}(1+|x|)^{a}
$$

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where the constants $M_{k}$ have sublinear growth:

$$
\begin{equation*}
\lim _{k \rightarrow \infty} \frac{M_{k}}{k}=\lim _{k \rightarrow \infty} \frac{M_{-k}}{k}=0 \tag{1.1}
\end{equation*}
$$

Then $\Delta f_{0}=-f_{0}$.
Conversely if $f$ is of polynomial growth $|f(x)| \leq M(1+|x|)^{a}$ and satisfies $\Delta f=$ $-f$, then $f_{k}=\Delta^{k} f=(-1)^{k} f$ satisfies $\left|f_{k}(x)\right| \leq M(1+|x|)^{a}$. The growth condition (1.1) is as weak as possible if polynomial growth of the functions is to be allowed. For example (in one dimension) if $f_{k}(x)=(-1)^{k}(x-2 k i) e^{i x}$ then $\frac{d^{2}}{d x^{2}} f_{k}=f_{k+1}$ and $\left|f_{k}(x)\right| \leq$ $(1+2|k|)(1+|x|)$ but $\frac{d^{2}}{d x^{2}} f_{0} \neq-f_{0}$. Even in the one dimensional case this gives the strengthening of Roe's theorem due to Burkill [2]: If $\left\langle f_{k}\right\rangle_{-\infty}^{\infty}$ satisfies $f_{k}^{\prime}(x)=f_{k+1}(x)$ and $\left|f_{k}(x)\right| \leq M(1+|x|)^{a}$, then $f_{0}(x)=A \sin (x+\alpha)$.

In Section 2 we extend Theorem 1 to any formally self-adjoint constant-coefficient differential operator. The proof has the flavor of Roe's original proof-using the growth conditions to show that the support of the Fourier transform of $f_{0}$ is contained in the unit sphere-but concluding that $f_{0}$ is an eigenfunction requires first showing that $f_{0}$ is a generalized eigenfunction of $\Delta$. A nonzero function $f$ is a generalized eigenfunction of $\Delta$ with eigenfunction $\lambda$ if and only if $(\Delta-\lambda)^{N} f=0$ for some $N \geq 1$. In one dimension the generalized eigenfunctions of $\Delta=\frac{d^{2}}{d x^{2}}$ were characterized in [3] and [5]. In $\mathbb{R}^{n}$ the result is

THEOREM 2. If in Theorem 1 the sublinear growth condition (1.1) is replaced by the subexponential growth condition

$$
\lim _{k \rightarrow \infty} \frac{M_{k}}{(1+\epsilon)^{k}}=\lim _{k \rightarrow \infty} \frac{M_{-k}}{(1+\epsilon)^{k}}=0
$$

for all $\epsilon>0$, then $f_{0}$ is a generalized eigenfunction of $\Delta$ with eigenvalue $\lambda=-1$.
We shall extend this to all formally self-adjoint constant-coefficient differential operators in Section 2 (see Theorem 4). Although the following is well known to experts, it seems to be interesting enough to record here.

COROLLARY. A smooth function of polynomial growth is a generalized eigenfunction of $\Delta$ with eigenfunction -1 if and only if the support of its Fourier transform is contained in the unit sphere $\{\xi:|\xi|=1\}$.

Finally, we note that for $n \geq 2$ there are many eigenfunctions of $\Delta$ than have greater than polynomial growth. For example (when $n=2$ )

$$
f(x, y)=e^{2 i x+\sqrt{3} y}
$$

satisfies $\Delta f=-f$ and has exponential growth. However, it seems unlikely that it is possible to characterise eigenfunctions of $\Delta$ in the class of functions of exponential growth: Let $\phi$ be any continuous function on $[-\delta, \delta]$ and set

$$
f_{k}(x)=\int_{-\delta}^{\delta}(\sin \theta+i \cos \theta)^{2 k} \epsilon^{(\sin \theta+i \cos \theta) x} \phi(\theta) d \theta
$$

Then $\Delta f_{k}=\frac{d^{2} f_{k}}{d x^{2}}=f_{k+1}$ and

$$
\left|f_{k}(x)\right| \leq \int_{-\delta}^{\delta} e^{x \sin \theta}|\phi(\theta)| d \theta \leq M e^{(\sin \delta)|x|}
$$

for all $x$. But $f_{0}$ is not an eigenfunction.
2. The general results. Let $x_{1}, \ldots, x_{n}$ be the usual coordinates in $\mathbb{R}^{n}$ and $i^{2}=-1$. Set

$$
D_{j}=\frac{1}{i} \frac{\partial}{\partial x_{j}}
$$

The factor of $\frac{1}{i}$ is included to make the operator $D_{j}$ formally self-adjoint. For a multiindex $\alpha=\left(\alpha_{1}, \ldots, \alpha_{n}\right)$ and $\xi=\left(\xi_{1}, \ldots, \xi_{n}\right) \in \mathbb{R}^{n}$, let $\xi^{\alpha}=\xi_{1}^{\alpha_{1}} \cdots \xi_{n}^{\alpha_{n}}$ and $D^{\alpha}=$ $D_{1}^{\alpha_{1}} \cdots D_{n}^{\alpha_{n}}$. Let

$$
P(\xi)=\sum_{\alpha} a_{\alpha} \xi^{\alpha}
$$

be a polynomial in $\xi$ and let

$$
\begin{equation*}
L=P(D)=\sum_{\alpha} a_{\alpha} D^{\alpha} \tag{2.1}
\end{equation*}
$$

be the corresponding constant-coeffcient linear partial differential operator. If $P$ is realvalued then $L$ will be formally self-adjoint. We now state our main result.

Theorem 3. Suppose $P(\xi)=\sum_{\alpha} a_{\alpha} \xi^{\alpha}$ is real-valued and $L=P(D)$. Let $\left\langle f_{k}\right\rangle_{-\infty}^{\infty}$ be a sequence of complex-valued functions on $\mathbb{R}^{n}$ so that

$$
L f_{k}=f_{k+1}
$$

and

$$
\begin{equation*}
\left|f_{k}(x)\right| \leq M_{k}(1+|x|)^{a}, \tag{2.2}
\end{equation*}
$$

where $\left\langle M_{k}\right\rangle_{-\infty}^{\infty}$ satisfies the sublinear growth condition

$$
\begin{equation*}
\lim _{k \rightarrow \pm \infty} \frac{M_{|k|}}{k}=0 \tag{2.3}
\end{equation*}
$$

Then $f=f_{+}+f_{-}$, where $L f_{+}=f_{+}$and $L f_{-}=-f_{-}$. If 1 (or -1 ) is not in the range of $P$ then $f_{+}=0\left(\right.$ or $\left.f_{-}=0\right)$.

If $P(\xi)=-|\xi|^{2}$ then $L=\Delta$. Then $f_{0}=f_{-}$which yields Theorem 1. For operators such as the d'Alembertian $\square=\frac{\partial^{2}}{\partial x_{1}^{2}}+\cdots+\frac{\partial^{2}}{\partial x_{p}^{2}}-\frac{\partial^{2}}{\partial x_{p+1}^{2}}-\cdots-\frac{\partial^{2}}{\partial x_{n}^{2}}$, we see that both $f_{+} \neq 0$ and $f_{-} \neq 0$ are possible. (cf. Theorem 3.1 of [6]).

The theorem applies to a class of operators more general than differential operators. In $\mathcal{S}\left(\mathbb{R}^{n}\right)$, the Fourier transform and its inverse are given by

$$
\hat{f}(\xi)=\frac{1}{(2 \pi)^{\frac{n}{2}}} \int e^{-i \xi \cdot x} f(x) d x
$$

and

$$
\begin{equation*}
\check{f}(x)=\frac{1}{(2 \pi)^{\frac{n}{2}}} \int e^{i x \cdot \xi} f(\xi) d \xi \tag{2.4}
\end{equation*}
$$

By duality these definitions extend to the space of tempered distributions; i.e., the dual space $\mathcal{S}\left(\mathbb{R}^{n}\right)$. Then, for the operator in (2.1), we have

$$
\begin{equation*}
\widehat{L f}(\xi)=P(\xi) \hat{f}(\xi) \tag{2.5}
\end{equation*}
$$

This can be used to define an operator $L$ on the space of tempered distributions even when $P(\xi)$ is not a polynomial. Such operators are called multiplier operators or translation-invariant pseudo-differential operators. For example convolution operators $L f=\phi * f$ are of this type. We note that for our result to hold it suffices that $P(\xi)$ be smooth and that (for each multi-index $\alpha$ ) there be numbers $C$ and $N$ with

$$
\left|D^{\alpha} P(\xi)\right| \leq C(1+|\xi|)^{N} .
$$

Then $L$ (defined by (2.5)) is a linear operator on the space of tempered distributions. Although many such functions exist, for example $P(\xi)=C e^{-|\xi|^{2}}$, the theorem is most interesting when $L$ is a differential operator.

To prove the theorem we first show that the support of $\hat{f}_{0}$ is contained in the set $\{\xi$ : $|P(\xi)|=1\}$. Formally from (2.5) and the Fourier inversion formula (2.4) we get

$$
f_{k}(x)=L^{k} f_{0}(x)=\frac{1}{(2 \pi)^{\frac{n}{2}}} \int e^{i x \cdot \xi} P(\xi)^{k} \hat{f}_{0}(\xi) d \xi
$$

If this is to stay bounded (as $k$ varies), the support of $\hat{f}_{0}$ must be contained in the set $\{\xi:|P(\xi)|=1\}$. More precisely, we have

Proposition. (A) If a functionf satisfies, for $k=0,1,2, \ldots$

$$
\left|L^{k} f(x)\right| \leq M_{k}(1+|x|)^{a},
$$

where the constants $M_{k}$ satisfy the subexponential growth condition

$$
\begin{equation*}
\lim _{k \rightarrow \infty} \frac{M_{k}}{(1+\epsilon)^{k}}=0 . \tag{2.6}
\end{equation*}
$$

for all $\epsilon>0$, then

$$
\operatorname{spt}(\hat{f}) \subseteq\{\xi:|P(\xi)| \leq 1\}
$$

(B) If $\left\langle f_{k}\right\rangle_{k=-\infty}^{0}$ is a sequence of functions with $L f_{k}=f_{k+1}$, for $k \leq-1$,

$$
\left|f_{k}(x)\right| \leq M_{k}(1+|x|)^{a},
$$

and

$$
\lim _{k \rightarrow \infty} \frac{M_{|k|}}{(1+\epsilon)^{k}}=0
$$

for all $\epsilon>0$, then

$$
\operatorname{spt}\left(\hat{f}_{0}\right) \subseteq\{\xi:|P(\xi)| \geq 1\} .
$$

Unlike in Theorem 3 the function $P$ can be complex valued and the proposition will still hold. This proposition is very closely related to the results of Section 3 of Garbardo's paper [3]. He shows this is related to the Paley-Wiener-Schwartz theorem.

LEMMA. Let $\Phi$ and $\phi$ be $C^{\infty}$-functions with $\phi$ compactly supported. Assume $|\Phi(\xi)| \leq r<1$, for all $\xi \in \operatorname{spt}(\phi)$. For any sequence $\left\langle M_{k}\right\rangle_{k=0}^{\infty}$ of constants satisfying the subexponential growth condition (2.6) and, for any multi-index $\alpha$,

$$
\lim _{k \rightarrow \infty} M_{k}\left\|D^{\alpha}\left(\phi \Phi^{k}\right)\right\|_{L_{2}}=0
$$

PROOF. By the product rule there are constants $C(\beta, \gamma) \geq 0$ so that

$$
\begin{aligned}
\left\|D^{\alpha}\left(\phi \Phi^{k}\right)\right\|_{L_{2}} & =\left\|\sum_{\beta+\gamma=\alpha} C(\beta, \gamma) D^{\gamma} \phi D^{\beta} \Phi^{k}\right\|_{L_{2}} \\
& \leq \sum_{\beta+\gamma=\alpha} C(\beta, \gamma)\left\|D^{\gamma} \phi D^{\beta} \Phi^{k}\right\|_{L_{2}}
\end{aligned}
$$

Since $\operatorname{spt}\left(D^{\alpha} \phi\right) \subseteq \operatorname{spt}(\phi)$, we may assume $|\Phi(\xi)| \leq r<1$ on the support of $\psi=D^{\gamma} \phi$. It is (thus) enough to show:

$$
\lim _{k \rightarrow \infty} M_{k}\left\|\psi D^{\beta}\left(\Phi^{k}\right)\right\|_{L_{2}}=0
$$

Assume $k>|\beta|:=\beta_{1}+\cdots+\beta_{n}$. Writing $\Phi^{k}=\Phi \cdots \Phi$ and using the product rule gives a sum with $k^{|\beta|}$ terms. Each term is a product of $k$ factors, at least $k-|\beta|$ of which are $\Phi$. The other factors are of the form $D^{\gamma} \Phi$, where $0 \leq|\gamma| \leq|\beta|$. Setting

$$
\Psi_{\beta}(x)=\left(\max _{|\gamma| \leq|\beta|}\left|D^{\gamma} \Phi(x)\right|\right)^{|\beta|}
$$

and using $|\Phi(x)| \leq r$ on $\operatorname{spt}(\psi)$, we get

$$
\begin{aligned}
\left|\psi D^{\beta}\left(\Phi^{k}\right)\right| & \leq k^{|\beta|}|\psi(x)||\Phi(x)|^{k-|\beta|} \Psi_{\beta}(x) \\
& \leq k^{|\beta|} r^{k-|\beta|}|\psi(x)| \Psi_{\beta}(x)
\end{aligned}
$$

Thus

$$
M_{k}\left\|\psi D^{\beta}(\Phi)^{k}\right\|_{L_{2}} \leq\left\|\psi \Psi_{\beta}\right\|_{L_{2}} k^{|\beta|} M_{k} r^{k-|\beta|}
$$

The growth condition (2.6) implies that the right-hand side goes to zero as $k \rightarrow \infty$.
To prove part (A) of the proposition, it suffices to show $\langle\hat{f}, \phi\rangle=0$ if $\phi \in C_{0}^{\infty}\left(\mathbb{R}^{n}\right)$ and $\operatorname{spt}(\phi) \cap\{\xi:|P(\xi)| \leq 1\}=\emptyset$. Since $\operatorname{spt}(\phi)$ is compact, there is some $r<1$ so that $\frac{1}{|P(\xi)|} \leq r$, for all $\xi \in \operatorname{spt}(\phi)$. Then

$$
\begin{aligned}
\langle\hat{f}, \phi\rangle & =\left\langle P^{k} \hat{f}, \frac{\phi}{P^{k}}\right\rangle \\
& =\left\langle\widehat{L^{k}} f, \frac{\phi}{P^{k}}\right\rangle \\
& =\left\langle L^{k} f,\left(\frac{\phi}{P^{k}}\right)^{\wedge}\right\rangle
\end{aligned}
$$

Choose an integer $m$ with $2 m \geq 2 a+n+1$. A calculation, using the hypothesis of the proposition and the Cauchy-Schwartz inequality, implies

$$
\begin{aligned}
|\langle\hat{f}, \phi\rangle| & \leq \int\left|L^{k} f(x)\right|\left|\left(\frac{\phi}{P^{k}}\right)^{\wedge}\right| d x \\
& \leq M_{k} \int \frac{(1+|x|)^{a}}{\left(1+|x|^{2}\right)^{\frac{m}{2}}}\left(1+|x|^{2}\right)^{\frac{m}{2}}\left|\left(\frac{\phi}{P^{k}}\right)^{\wedge}\right| d x \\
& \leq M_{k}\left(\int \frac{(1+|x|)^{2 a}}{\left(1+|x|^{2}\right)^{m}} d x\right)^{\frac{1}{2}}\left(\int\left(1+|x|^{2}\right)^{m}\left|\left(\frac{\phi}{P^{k}}\right)^{\wedge}\right|^{2} d x\right)^{\frac{1}{2}} \\
& =M_{k} C_{1}(a, m, n)\left(\int\left(1+|x|^{2}\right)^{m}\left|\left(\frac{\phi}{P^{k}}\right)^{\wedge}\right|^{2} d x\right)^{\frac{1}{2}} .
\end{aligned}
$$

By a standard estimate (cf. [4], Chapter 1), there is a constant $C_{2}(m, n)$ with

$$
\left(\int\left(1+|x|^{2}\right)^{m}|\hat{f}(x)|^{2} d x\right)^{\frac{1}{2}} \leq C_{2} \sum_{|\alpha| \leq m}\left\|D^{\alpha} f\right\|_{L_{2}}
$$

Using this in the above leads to

$$
|\langle\hat{f}, \phi\rangle| \leq C_{3}(m, n, a) M_{k} \sum_{|\alpha| \leq m}\left\|D^{\alpha}\left(\frac{\phi}{P^{k}}\right)\right\|_{L_{2}} .
$$

By the lemma the right-hand side of this goes to zero as $k \rightarrow \infty$, and so $\langle\hat{f}, \phi\rangle=$ 0 . This completes the proof of part (A); part (B) is similar. Let $\phi \in C_{0}^{\infty}\left(\mathbb{R}^{n}\right)$ so that $\operatorname{spt}(\phi) \cap\{\xi:|P(\xi)| \geq 1\}=\emptyset$. We shall show that $\left\langle\hat{f}_{0}, \phi\right\rangle=0$. Then, for some $r<1$, the inequality $|P(\xi)| \leq r$ holds for all $\xi$ in $\operatorname{spt}(\phi)$. Thus

$$
\begin{aligned}
\left\langle\hat{f_{0}}, \phi\right\rangle & =\left\langle\widehat{L^{k} f_{-k}}, \phi\right\rangle \\
& =\left\langle P^{k} \widehat{f_{-k}}, \phi\right\rangle \\
& =\left\langle f_{-k}, \widehat{P^{k}} \phi\right\rangle
\end{aligned}
$$

The rest follows as in part (A).
We now prove Theorem 3. We first assume that -1 is not a value of $P(\xi)$, and show that $L f_{0}=f_{0}$. Let $S=\{\xi: P(\xi)=1\}$. That $\operatorname{spt}\left(\hat{f}_{0}\right) \subseteq S$ follows from the growth conditions on the sequence $\left\langle f_{k}\right\rangle_{-\infty}^{\infty}$, the proposition, and the assumption that $P(\xi) \neq-1$.

The topology on the space $\mathcal{S}\left(\mathbb{R}^{n}\right)$ is defined by the seminorms

$$
\|\phi\|_{N, m}=\sup _{x} \sum_{|\alpha| \leq N}(1+|x|)^{m}\left|D^{\alpha} \phi(x)\right| .
$$

Therefore, since $\hat{f}_{0}$ is a continuous linear functional on $\mathcal{S}\left(\mathbb{R}^{n}\right)$, there is a constant $C$ and integers $m$ and $N$ so that

$$
\begin{equation*}
\left|\left\langle\hat{f}_{0}, \phi\right\rangle\right| \leq C\|\phi\|_{N, m}, \tag{2.7}
\end{equation*}
$$

for all $\phi \in S\left(\mathbb{R}^{n}\right)$. Therefore (as a distribution) $\hat{f}_{0}$ is of order $\leq N$.

Claim. For this $N$,

$$
\begin{equation*}
(P-1)^{N+1} \hat{f}_{0}=0 \tag{2.8}
\end{equation*}
$$

To simplify notation set $h:=(P-1)$. Then we need to show, for any compactly supported $C^{\infty}$ function $\phi$, that

$$
\left\langle h^{N+1} \hat{f}_{0}, \phi\right\rangle:=\left\langle\hat{f}_{0}, h^{N+1} \phi\right\rangle=0
$$

Let $g: \mathbb{R} \rightarrow[0,1]$ be a $C^{\infty}$ function with $g=1$ on $[-1 / 2,1 / 2]$ and $g=0$ outside $(-1,1)$. Set

$$
g_{r}(t):=g\left(\frac{t}{r}\right)
$$

Letting $B=\max \left\{\left|g^{(k)}(t)\right|: t \in[-1,1], k \leq N\right\}$, we have

$$
\begin{equation*}
\left|g_{r}^{(k)}(t)\right| \leq \frac{B}{r^{k}} \leq \frac{B}{|t|^{k}}, \tag{2.9}
\end{equation*}
$$

for all $k \leq N$. Set

$$
H_{r}=g_{r}(h) h^{N+1} \phi .
$$

Then $H_{r}=h^{N+1} \phi$ in a neighborhood of $\{\xi: h(\xi)=0\}=\{\xi: P(\xi)=1\} \supseteq \operatorname{spt} \hat{f}_{0}$. Thus by (2.7) we have

$$
\left|\left\langle\hat{f}_{0}, h^{N+1} \phi\right\rangle\right|=\left|\left\langle\hat{f}_{0}, H_{r}\right\rangle\right| \leq C\left\|H_{r}\right\|_{N, m} .
$$

To verify (2.8), it suffices to demonstrate $\left\|H_{r}\right\|_{N, m} \rightarrow 0$ as $r \rightarrow 0$. Write $D^{k}$ for any $D^{\alpha}$ with $|\alpha|=k$ and $(D h)^{k}$ for a product $D_{j_{1}} h \cdots D_{j_{k}} h$ of $k$ first order partial derivatives of $h$. Then (ignoring factors of $i$ ) and assuming $k \leq N$

$$
D^{k}\left(H_{r}\right)=D^{k}\left(g_{r}(h) h^{N+1} \phi\right)
$$

is a sum of terms of the form

$$
T_{r}\left(k_{1}, \ldots, k_{l+1}\right)=g_{r}^{\left(k_{1}\right)}(h)(D h)^{k_{1}}\left(D^{k_{2}} h\right) \cdots\left(D^{k_{l}} h\right) h^{N+1-k_{2}-\cdots-k_{l}} D^{k_{l+1}} \phi
$$

where $k_{1}+\cdots+k_{l+1}=k$. Since the support of $\phi$ is compact, there is a bound $K$ so that, for all $x \in \operatorname{spt}(\phi)$,

$$
(1+|x|)^{m}\left|D^{\alpha} \phi(x)\right| \leq K \text { and }\left|D^{\alpha} h(x)\right| \leq K,
$$

whenever $|\alpha| \leq N$. Since $|h| \leq r$ on the support of $H_{r}$, the inequality (2.9) implies that on the support of $T_{r}\left(k_{1}, \ldots, k_{l+1}\right)$,

$$
\begin{aligned}
(1+|x|)^{m}\left|T_{r}\left(k_{1}, \ldots, k_{l+1}\right)\right| & \leq \frac{B}{|h|^{k_{1}}}|h|^{N+1-k_{2}-\cdots-k_{l}} K K^{k_{l-1}} K^{k_{1}} \\
& =B|h|^{N+1-k+k_{l+1}} K^{l+k_{1}} \\
& \leq B K^{l+k_{1}} r^{N+1-k+k_{l+1}} .
\end{aligned}
$$

But $k \leq N$ so this goes to zero as $r \rightarrow 0$. The sum defining $\left\|H_{r}\right\|_{N, m}$ is a finite sum of terms of this type and so $\left\|H_{r}\right\|_{N, m} \rightarrow 0$ as $r \rightarrow 0$. This completes the proof of the claim.

Inverting the Fourier transform (2.8) yields that

$$
\begin{equation*}
(L-1)^{N+1} f_{0}=0 . \tag{2.10}
\end{equation*}
$$

This equation implies

$$
\operatorname{span}\left\{f_{0}, f_{1}, f_{2}, \ldots\right\}=\operatorname{span}\left\{f_{0}, L f_{0}, L^{2} f_{0}, \ldots\right\}=\operatorname{span}\left\{f_{0}, \ldots, L^{N} f_{0}\right\}
$$

We shall now show that we can take $N=0$ in (2.10). If not then $(L-1) f_{0} \neq 0$. Let $K$ be the largest positive integer so that $(L-1)^{K} f \neq 0$. Clearly $K \leq N$. Thus

$$
f:=(L-1)^{K-1} f_{0} \in \operatorname{span}\left\{f_{0}, \ldots, f_{N}\right\}
$$

will satisfy

$$
\begin{equation*}
(L-1)^{2} f=0 \text { and }(L-1) f \neq 0 . \tag{2.11}
\end{equation*}
$$

Write

$$
f=a_{0} f_{0}+\cdots+a_{N} f_{N},
$$

for constants $a_{0}, \ldots, a_{N}$. Then

$$
L^{k} f=a_{0} f_{k}+\cdots+a_{N} f_{N+k} .
$$

If $C_{k}=\left|a_{0}\right| M_{k}+\cdots+\left|a_{N}\right| M_{k+N}$, then this and (2.2) imply

$$
\begin{equation*}
\left|\left(L^{k} f\right)(x)\right| \leq C_{k}(1+|x|)^{a} . \tag{2.12}
\end{equation*}
$$

By (2.3) these satisfy the sublinear growth condition

$$
\begin{equation*}
\lim _{k \rightarrow \infty} \frac{C_{k}}{k}=0 \tag{2.13}
\end{equation*}
$$

An induction using (2.11) implies for $k \geq 2$ that

$$
\begin{aligned}
L^{k} f & =k L f-(k-1) f=k(L-1) f+f . \\
|((L-1) f)(x)| & \leq \frac{1}{k}\left|\left(L^{k} f\right)(x)\right|+\frac{|f(x)|}{k} \leq \frac{C_{k}}{k}(1+|x|)^{a}+\frac{|f(x)|}{k}
\end{aligned}
$$

Letting $k \rightarrow \infty$ and using (2.13) implies $(L-1) f=0$. But this contradicts (2.11). Consequently, $N=0$ in (2.10). This completes the proof in the case that -1 is not in the range of $P$.

In the case that +1 is not in the range of $P$ we apply the same argument to $-L$ to conclude $L f_{0}=-f_{0}$. In the general case, let $L_{0}=L^{2}$. Then $\widehat{L_{0} f}(\xi)=P(\xi)^{2} \hat{f}(\xi) . L_{0} f_{2 k}=$ $f_{2(k+1)}$ and $P(\xi)^{2} \neq-1$. Thus we can (as before) conclude, for the sequence $\left\langle f_{2 k}\right\rangle_{k=-\infty}^{\infty}$ that

$$
L_{0} f_{0}=L^{2} f_{0}=f_{0}
$$

Set $f_{+}=\frac{1}{2}\left(f_{0}+L f_{0}\right)$ and $f_{-}=\frac{1}{2}\left(f_{0}-L f_{0}\right)$. Then $f=f_{+}+f_{-}, L f_{+}=f_{+}$, and $L f_{-}=-f_{-}$. This completes the proof of Theorem 3.

THEOREM 4. If, in Theorem 3, we replace (2.3) with

$$
\begin{equation*}
\lim _{k \rightarrow \pm \infty} \frac{M_{|k|}}{(1+\epsilon)^{k}}=0 \tag{2.14}
\end{equation*}
$$

for all $k>0$, then the span of $\left\langle f_{k}\right\rangle$ is finite dimensional. Moreover, $f_{0}=f_{+}+f_{-}$, where, for some integer $N,(L-1)^{N} f_{+}=0$ and $(L+1)^{N} f_{-}=0$. Thus $f_{+}$(or $\left.f_{-}\right)$is a generalized eigenfunction of $L$ with eigenvalue +1 (or -1 ).

The proof will be based on the following result from linear algebra.
LEMMA. Let $X$ be a finite dimensional complex vector space, and let $A: X \rightarrow X$ be a linear map with eigenvalues $\lambda_{1}, \ldots, \lambda_{p}$. Then $X=X_{1} \oplus \cdots \oplus X_{p}$, where $X_{j}=$ $\operatorname{ker}\left(\left(A-\lambda_{j}\right)^{N}\right)$ and $N=\operatorname{dim} X$.

This can be deduced from the Jordan normal form. (cf. [1], Chapter 10.)
We first prove Theorem 4 under the assumption that $P(\xi) \neq-1$. Using the growth condition (2.14) and the proposition, we may still conclude that $\operatorname{spt}\left(\hat{f}_{0}\right) \subseteq S=\{\xi$ : $P(\xi)=1\}$. But then, as before, we can conclude that (2.10) holds. But this is enough to complete the proof in this case. A similar argument shows that if $P(\xi) \neq 1$, then $(L+1)^{N} f_{0}=0$.

In the general case we again let $L_{0}=L^{2}$ and $P_{0}=P^{2}$. Then $P_{0}(\xi) \neq-1$ and the span of $\left\langle f_{2 k}\right\rangle$ is finite dimensional. The map $L$ takes the span of $\left\langle f_{2 k}\right\rangle$ onto the span of $\left\langle f_{2 k+1}\right\rangle$. Thus $X$ is finite dimensional. Any $f \in X$ will have $\operatorname{spt}(f)$ inside the set defined by $P(\xi)= \pm 1$. From this it is not hard to show the only possible eigenvalues of $L$ restricted to $X$ are +1 and -1 . The result now follows from the last lemma.

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