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## Anne Berthier <br> On the point spectrum of Schrödinger operators

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# ON THE POINT SPECTRUM OF SCHRÖDINGER OPERATORS 

By Anne BERTHIER

## 1. Introduction

This paper is an extension of a work [2] on the spectral analysis of partial differential operators of Schrödinger type. The problem was the following: Let A be a compact subset of $\mathbb{R}^{n}, \Sigma$ a finite interval in $\mathbb{R}$ and H a self-adjoint elliptic differential operator in the complex Hilbert space $\mathscr{H}=\mathrm{L}^{2}\left(\mathbb{R}^{n}\right)$. We define $\mathrm{F}(\Sigma)$ to be the spectral projection of H associated with the interval $\Sigma$ and $\mathrm{E}(\mathrm{A})$ the multiplication operator by the characteristic function $\chi_{\mathrm{A}}$ of A . Do there exist vectors in $\mathrm{L}^{2}\left(\mathbb{R}^{n}\right)$ which are contained both in the range $\mathrm{E}(\mathrm{A}) \mathscr{H}$ of $\mathrm{E}(\mathrm{A})$ and in $\mathrm{F}(\Sigma) \mathscr{H}$ ?

It turns out that the closed subspace $\mathscr{H}_{p}(\mathrm{H})$ generated by the set of eigenvectors of H plays a different role from the subspace $\mathscr{H}_{c}(\mathrm{H})=\mathscr{H}_{p}(\mathrm{H})^{\perp}$ associated with the continuous spectrum of H . Notice that it is shown in [2], under regularity and integrability conditions on the coefficients of the differential operator, that there do not exist vectors of $\mathscr{H}_{c}(\mathrm{H})$ which belong both to $\mathrm{E}(\mathrm{A}) \mathscr{H}$ and to $\mathrm{F}(\Sigma) \mathscr{H}$. On the other hand, to prove the non-existence of vectors in $\mathscr{H}_{p}(\mathrm{H})$ belonging to $\mathrm{E}(\mathrm{A}) \mathscr{H} \cap \mathrm{F}(\Sigma) \mathscr{H}$, we used an unique continuation theorem for solutions of the differential equation associated with H . Now, if for example $\mathrm{H}=-\Delta+\mathrm{V}$, where V is the multiplication operator by a real function $v(\vec{x})$, the known results on unique continuation require a condition $\mathrm{L}^{\infty}\left(\mathbb{R}^{n} \backslash \mathrm{~N}\right)$ on $v$, where N is a closed set of measure zero such that $\mathbb{R}^{n} \backslash \mathrm{~N}$ is connected ([3], [5]).

In the present paper, we propose to show that:

$$
\begin{equation*}
\mathscr{H}_{p}(\mathrm{H}) \cap \mathrm{E}(\mathrm{~A}) \mathscr{H} \cap \mathrm{F}(\Sigma) \mathscr{H}=\{0\} \tag{1}
\end{equation*}
$$

by imposing only an integrability condition on the function $v$. More precisely, we will prove (1) under the hypothesis that $v \in \mathrm{~L}_{\mathrm{Loc}}^{s}\left(\mathbb{R}^{n}\right)$ with $s=2$ if $n=1,2,3$ and $s>n-2$ if $n \geqq 4$.

This result shows that, under the above conditions on $v$, the operator $-\Delta+v$ has no eigenvector with compact support. This is essentially the content of our Theorem 1 in paragraph 2. (In the case $n=1$, one obtains ordinary differential operators for which results of this type have been known for a long time [9]).

[^0]This result is also interesting from the point of view of "non-existence of positive eigenvalues of the operator H". In the literature (for example [2], [12]) the non-existence of positive eigenvalues is obtained in two steps:
(i) under suitable decay conditions at infinity on the function $v$, it is shown that all eigenfunctions $f$ associated with a strictly positive eigenvalue of H have compact support;
(ii) then one imposes suitable local conditions on $v\left(\right.$ e.g. $v \in \mathrm{~L}_{\text {Loc }}^{\infty}\left(\mathbb{P}^{n} \backslash \mathrm{~N}\right)$ in order to apply the unique continuation theorem, which then leads to $f \equiv 0$. It turns out that the nonexistence of positive eigenvalues is also obtained by assuming in (ii) as a local condition that $v \in \mathrm{~L}_{\text {Loc }}^{s}\left(\mathbb{R}^{n}\right)$ with $s=2$ if $n=1,2,3$ and $s>n-2$ if $n \geqq 4$ (Thm. 2).

Finally our method implies also the spectral continuity of a class of Schrödinger operators with periodic potentials $v(\vec{x})$.

The organization of the paper is a follows: first we give the principal results and deduce Theorems 1 and 2 from Theorem 3 in section 2, and we introduce a direct integral representation of Schrödinger operators in section 3. This representation will be used in section 4 for proving Theorem 3. The principal estimate of the proof is the subject of the last section 5 .

## 2. Statements of the results

Let $v: \mathbb{R}^{n} \rightarrow \mathbb{R}$ be a measurable function. We always suppose that:

$$
\begin{equation*}
v \in \mathrm{~L}_{\mathrm{Loc}}^{s}\left(\mathbb{R}^{n}\right) \quad \text { with } s=2 \text { if } n=1,2,3 ; \quad s>n-2 \quad \text { if } n \geqq 4 . \tag{2}
\end{equation*}
$$

Notice that $s>n-2$ in all cases.
The function $v$ will be called periodic if there exist $n$ linearly independent vectors $\vec{a}_{1}, \ldots, \vec{a}_{n} \in \mathbb{R}^{n}$ such that $v\left(\vec{x}+\vec{a}_{i}\right)=v(\vec{x})$ for all $\vec{x} \in \mathbb{R}^{n}$. A periodic function will be called ortho-periodic if:

$$
\begin{equation*}
\vec{a}_{j} \cdot \vec{a}_{k}=\mathrm{L}^{2} \delta_{j k} \tag{3}
\end{equation*}
$$

with $\mathrm{L}>0$, i.e. if the vectors of the form $\sum_{i=1}^{n} \alpha_{i} \cdot \vec{a}_{i}, 0 \leqq \alpha_{i}<1$, define a cube $\mathrm{C}^{n}$ with side L .
We denote by $\hat{\mathrm{H}}$ the symmetric operator:

$$
\begin{equation*}
\hat{\mathrm{H}}=-\Delta+v(\vec{x}) \tag{4}
\end{equation*}
$$

with domain $\mathrm{D}(\hat{\mathrm{H}})=\mathrm{C}_{0}^{\infty}\left(\mathbb{R}^{n}\right)$ and by $\mathrm{H}_{0}$ the unique self-adjoint extension of $\hat{\mathrm{H}}_{0}=-\Delta$, $\mathrm{D}\left(\hat{\mathrm{H}}_{0}\right)=\mathrm{C}_{0}^{\infty}\left(\mathbb{R}^{n}\right)$. Let H a self-adjoint extension of $\hat{H}$. We have the following lemma:

Lemma 1. - Assume that (2) and one of the following conditions are satisfied:
(i) $v$ is periodic;
(ii) $v \in \mathrm{~L}^{\infty}\left(\left\lceil\mathrm{B}_{\mathrm{R}}\right)\right.$ where $\mathrm{B}_{\mathrm{R}}=\left\{\vec{x} \in \mathbb{R}^{n}| | \vec{x} \mid \leqq \mathrm{R}\right\}$ and $\left\lceil\mathrm{B}_{\mathrm{R}}\right.$ denotes the complement of $\mathrm{B}_{\mathrm{R}}$.

Then:
(a) $v$ is $\mathrm{H}_{0}$-bounded with $\mathrm{H}_{0}$-bound 0 ;

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(b) $\hat{\mathrm{H}}$ is essentially self-adjoint;
(c) $\mathrm{D}(\mathrm{H})=\mathrm{D}\left(\mathrm{H}_{0}\right)$, where H is the unique self-adjoint extension of H .

Proof. - (b) and (c) follow from (a) by using the Kato-Rellich Theorem ([7], Chapt. 5.4.1). Under hypothesis (i), (a) follows from Theorem XIII. 96 of [11], whereas under the assumption (ii), (a) can be proved by the method used in the proof of Lemma 3 in [10]. Both cases are treated in [4].

We now state our principal results. In Theorem 2 we choose as conditions on the potential $v$ at infinity those used in [4].

Theorem 1. - Let $v \in \mathrm{~L}_{\mathrm{L} . \text { sc }}^{s}\left(\mathbb{R}^{n}\right)$ uith s satisfying (2) and let H be a self-adjoint extension of H : (a) suppose that $f \in \mathrm{~L}^{2}\left(\mathbb{R}^{n}\right)$ satisfies $\mathrm{H} f=\lambda f$ for some $\lambda \in \mathbb{R}$ and $\mathrm{E}(\mathrm{A}) f=f$ for some compact subset A of $\mathbb{R}^{n}$. (i.e. $f$ is an eigenvector of H with compact support in $\mathbb{R}^{n}$ ). Then $f=0$;
(b) for each compact subset A of $\mathbb{R}^{n}$ and each bounded interval $\Sigma$, one has:

$$
\mathscr{H}_{p}(\mathrm{H}) \cap \mathrm{E}(\mathrm{~A}) \mathscr{H} \cap \mathrm{F}(\Sigma) \mathscr{H}=\{0\} .
$$

Theorem 2. - Suppose that:
(i) $v \in \mathrm{~L}^{s}\left(\mathrm{~B}_{\mathrm{R}}\right)$ with $s$ satisfying (2) for some $\mathrm{R}<\infty$;
(ii) $v=v_{1}+v_{2}$ such that:
$(\alpha) v_{1}, v_{2} \in \mathrm{~L}^{\infty}\left(\left\lceil\mathrm{B}_{\mathrm{R}}\right)\right.$,
$(\beta)|\vec{x}| v_{1}(\vec{x}) \rightarrow 0 \quad$ as $\quad|\vec{x}| \rightarrow \infty$,
$(\gamma) v_{2}(\vec{x}) \rightarrow 0 \quad$ as $\quad|\vec{x}| \rightarrow \infty$,
$(\delta) r \mapsto v_{2}(r,$.
is differentiable as a function from $(\mathrm{R}, \infty)$ to $\mathrm{L}^{\infty}\left(\mathrm{S}^{n-1}\right)$, and $\lim \sup \partial v_{2} / \partial r \leqq 0$. ( $\mathrm{S}^{n-1}$ denotes the unit sphere in $\mathbb{R}^{n}$.)

Then $\mathrm{H}=\mathrm{H}_{0}+\mathrm{V}$ has no eigenvalues in $(0, \infty)$.
Theorem 3. - Let $v$ be ortho-periodic and $v \in \mathrm{~L}_{\mathrm{Loc}}^{s}\left(\mathbb{R}^{n}\right)$ with $s$ satisfying (2). Then the spectrum of $\mathrm{H}=\mathrm{H}_{0}+\mathrm{V}$ is purely continuous.

Remark 1. - By following the proof of Theorem XIII. 100 in [11], it is possible to show that the operator H in Theorem 3 is absolutely continuous. Other comments on Theorem 3 will be made at the end of this paper.

Remark 2. - Contrarily to [2], where the operator $\hat{\mathrm{H}}$ was defined by:

$$
\hat{\mathrm{H}}=\sum_{j, k=1}^{n} a_{j k}\left(-i \frac{\partial}{\partial x_{j}}+b_{j}(\vec{x})\right)\left(-i \frac{\partial}{\partial x_{k}}+b_{k}(\vec{x})\right)+\mathrm{V}(\vec{x}),
$$

we assume here that the vector potential $\vec{b}=\left\{b_{k}\right\}$ is equal to zero. It is possible to generalize Theorem 1 to the case where $\vec{b} \neq 0$.

Theorem 2 follows from results of [11] and [6], and from Theorem 1 as indicated in the introduction. (If $\mathrm{H} f=\lambda f$ with $\lambda>0$, then $f$ has compact support by Theorem XIII. 58 of
[11], and consequently $f=0$ by our Theorem 1.) Theorem $1(a)$ is deduced from Theorem 3: By the proof of Proposition 4 of [2], the vector $f$ belongs to $\mathrm{D}\left(\mathrm{H}_{0}\right) \cap \mathrm{D}(\mathrm{V})$ and $\mathrm{H} f=\mathrm{H}_{0} f+\mathrm{VE}(\mathrm{A}) f$. Let $w$ be an ortho-periodic function such that $w \in \mathrm{~L}_{\mathrm{Loc}}^{s}\left(\mathbb{R}^{n}\right)$ and $w(\vec{x})=v(\vec{x})$ for $\vec{x} \in \mathrm{~A}$. If $\mathrm{H}_{1}$ denotes the periodic Schrödinger operator $\mathrm{H}_{1}=\mathrm{H}_{0}+\mathrm{W}$ then $\mathrm{H}_{1} f=\mathrm{H} f=\lambda f$. Therefore we deduce from Theorem 3 that $f=0$.

To show Theorem $1(b)$, let $\mathrm{S}=\mathrm{E}(\mathrm{A}) \cap \mathrm{F}(\Sigma)$ (the orthogonal projection with range $\mathrm{E}(\mathrm{A}) \mathscr{H} \cap \mathrm{F}(\Sigma) \mathscr{H})$ and suppose that $f \in \mathscr{H}_{p}(\mathrm{H})$ satisfies $\mathrm{S} f=f . \quad f$ is a linear combination of eigenvectors of H , i.e. $f=\sum_{k} \alpha_{k} . g_{k}$, where $\mathrm{H} g_{k}=\lambda_{k} g_{k}$ with $\lambda_{k} \in \Sigma$. It follows that:

$$
\mathbf{S} f=f=\sum_{k} \alpha_{k} \mathrm{~S} g_{k}
$$

Now, by Proposition 2 of [2], S commutes with H ; in particular $\mathrm{HS} g_{k}=\mathrm{SH} g_{k}=\lambda_{k} \mathrm{~S} g_{k}$. This implies that each $\mathrm{S} g_{k}$ is an eigenvector of H of compact support in A, hence $\mathrm{S} g_{k}=0$ by the part $(a)$ of Theorem 1. We deduce from this that $f=\sum_{k} \alpha_{k} \mathrm{~S} g_{k}=0$. The condition " $\Sigma$ bounded" is fundamental: we can choose a potential V such that $\mathscr{H}_{p}(\mathrm{H})=\mathscr{H}$, i. e. such that the eigenvectors of $\mathscr{H}$ generate $\mathscr{H}$. In this case, we have:

$$
\left.\mathscr{H}_{p}(\mathrm{H}) \cap \mathrm{E}(\mathrm{~A}) \mathscr{H}=\mathrm{E}(\mathrm{~A}) \mathscr{H} \neq\{0\} .\right]
$$

## 3. Reduction of the translation group of the lattice

In this part, let $v$ be an ortho-periodic potential. In a natural way, this implies a decomposition of the Hilbert space $\mathscr{H}=\mathrm{L}^{2}\left(\mathbb{R}^{n}\right)$ and of the operators H and $\mathrm{H}_{0}$ into direct integrals. This decomposition will be used in the next part for the proof of Theorem 3.

The potential $v$ satisfies $v\left(\vec{x}+\vec{a}_{i}\right)=v(\vec{x})$ where $\vec{a}_{1}, \ldots, \vec{a}_{n}$ are as in (3). The points of the form $\vec{z}=\sum_{i=1}^{n} q_{i} \vec{a}_{i}, \vec{q}=\left\{q_{i}\right\} \in \mathbb{Z}^{n}$, form a cubic lattice in $\mathbb{R}^{n}$ which is invariant under the translations:

$$
\vec{z} \mapsto \vec{z}+\sum_{i} q_{i}^{\prime} \vec{a}_{i}, \quad \overrightarrow{q^{\prime}} \in \mathbb{Z}^{n}
$$

In $L^{2}\left(\mathbb{R}^{n}\right)$. we consider the unitary representation $U(\vec{q})$ of the additive group $\mathbb{Z}^{n}$ given by:

$$
\begin{equation*}
[\mathrm{U}(\vec{q}) f](\vec{x})=f\left(\vec{x}-\sum_{i} q_{i} \vec{a}_{i}\right)=f(x-\mathbf{L} \vec{q}) \tag{5}
\end{equation*}
$$

where we have written $\sum_{i} q_{i} \vec{a}_{i}=\mathbf{L} \vec{q}$, assuming that the directions of the $\vec{a}_{i}$ coïncide with Cartesian coordinate system.

We also introduce the reciprocal lattice which is the set of points of the following form:

$$
\vec{z}=\sum_{i=1}^{n} q_{i} \vec{e}_{i}, \quad \vec{q} \in \mathbb{Z}^{n}
$$

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where the vectors $\vec{e}_{1}, \ldots, \vec{e}_{n}$ are defined by:

$$
\begin{equation*}
\vec{e}_{i} \cdot \vec{a}_{k}=2 \pi \delta_{i k} \tag{6}
\end{equation*}
$$

We may write $\vec{z}=\mathrm{E} \vec{q}$, with $\mathrm{E}=2 \pi \mathrm{~L}^{-1}$. Let again:

$$
\Gamma^{n}=\left\{k \in \mathbb{R}^{n} \mid k=\sum_{i=1}^{n} \lambda_{i} e_{i}, 0 \leqq \lambda_{i}<1\right\}
$$

Consider the Hilbert space $\mathscr{G}$ of square-integrable functions $f: \Gamma^{n} \rightarrow l_{n}^{2} \equiv l^{2}\left(\mathbb{Z}^{n}\right)$ :

$$
\mathscr{G}=\mathrm{L}^{2}\left(\Gamma^{n} ; l_{n}^{2}\right)
$$

We write $f(\vec{k})_{\vec{q}}$ for the component $\vec{q}\left(\vec{q} \in \mathbb{Z}^{n}\right)$ of $f$ at the point $\vec{k} \in \mathbb{Z}^{n}$. Thus, we have:

$$
\|f\|_{\mathscr{\mathcal { G }}}^{2}=\int_{\Gamma^{n}} d k \sum_{\overrightarrow{\bar{u}} \in, / n}\left|f(\vec{k})_{\vec{q}}\right|^{2}
$$

Now, let $\mathscr{U}: \mathscr{H} \rightarrow \mathscr{G}$ be the operator defined by:

$$
\begin{equation*}
(\mathscr{U} f)(\vec{k})_{\vec{q}}=\hat{f}(\vec{k}+\mathrm{E} \vec{q}), \tag{7}
\end{equation*}
$$

where $\hat{f}$ is the Fourier transform of the function $f$ :

$$
\hat{f}(\xi)=(2 \pi)^{-n / 2} \operatorname{Lim} \int_{\mathbb{R}^{n}} d x \exp (-i \vec{x} \cdot \vec{\xi}) f(\vec{x})
$$

It follows from Plancherel's Theorem that the operator $\mathscr{U}$ is unitary, and its inverse is given by:

$$
\mathscr{\mathscr { F }}\left[\mathscr{U}^{-1}\{f(.)\}\right](\xi)=f(\vec{k})_{\dot{q}},
$$

where $\vec{q} \in \mathbb{Z}^{n}$ and $\vec{k} \in \Gamma^{n}$ are determined by $\vec{k}+\mathrm{E} \vec{q}=\vec{\xi}$. If $\vec{m} \in \mathbb{Z}^{n}$, one has:

$$
\begin{equation*}
[\mathscr{U} \mathrm{U}(\vec{m}) f](\vec{k})_{\vec{q}}=\exp (-i \mathrm{~L} \vec{k} \cdot \vec{m})(\mathscr{U} f)(\vec{k})_{\vec{q}} \tag{8}
\end{equation*}
$$

i. e. $\mathscr{U} \mathrm{U}(m) \mathscr{U}^{-1}$ is diagonalizable in $\mathscr{G}$ (i.e. a multiplication operator by a function of $\vec{k})$. As the functions $\{\exp (\mathrm{i} L \vec{k} \cdot \vec{m})\}_{\vec{m} \in \mathbb{Z}^{n}}$ form a basis of $\mathrm{L}^{2}\left(\Gamma^{n}\right)$, each bounded diagonalizable operator is a function of $\left\{\mathscr{U} \mathrm{U}(m) \mathscr{U}^{-1}\right\}$. As $\mathrm{H}_{0}, \mathrm{~V}$ and H commute with every $\mathrm{U}(\vec{m})$, these operators commute with each diagonalizable operator, i.e. $\mathscr{U} \mathrm{H}_{0} \mathbb{U}^{-1}$, $\mathscr{U} \mathrm{V} \mathscr{U}^{-1}$ and $\mathscr{U} \mathrm{H} \mathscr{U}^{-1}$ are decomposable in $\mathrm{L}^{2}\left(\Gamma^{n} ; l_{n}^{2}\right)$. Therefore there exist in $l_{n}^{2}$ measurable families of self-adjoint operators $\mathrm{H}_{0}(\vec{k}), \mathrm{V}(\vec{k})$ and $\mathrm{H}(\vec{k})\left(\vec{k} \in \Gamma^{n}\right)$ such that, for $f \in \mathrm{D}\left(\mathrm{H}_{0}\right)$ :

$$
\left\{\begin{array}{c}
\left(\mathscr{U} \mathrm{H}_{0} f\right)(\vec{k})=\mathrm{H}_{0}(\vec{k}) f(\vec{k}),  \tag{9}\\
(\mathscr{U} \mathrm{V} f)(\vec{k})=\mathrm{V}(\vec{k}) f(\vec{k}), \\
(\mathscr{U} \mathrm{H} f)(\vec{k})=\mathrm{H}(\vec{k}) f(\vec{k})
\end{array}\right.
$$

Now let us give the explicit form and the properties of these three families of operators.

Lemma 2. - (i) $\mathrm{H}_{0}(\vec{k})$ is the self-adjoint multiplication operator in $l_{n}^{2}$ by $\varphi_{\vec{k}}(\vec{q})=(\vec{k}+\mathrm{E} \vec{q})^{2}:$ If $g=\left\{g_{\vec{q}}\right\} \in l_{n}^{2}$, then:

$$
\left(\mathrm{H}_{0}(\vec{k}) g\right)_{\vec{q}}=(\vec{k}+\mathrm{E} \vec{q})^{2} g_{\vec{q}}
$$

(ii) the domain of $\mathrm{D}\left(\mathrm{H}_{0}(\vec{k})\right)$ is independent of $\vec{k}$ and is given by:

$$
\mathrm{D}\left(\mathrm{H}_{0}(\vec{k})\right)=\mathrm{D}_{0}=\left\{\left.g \in l_{n}^{2}\left|\sum_{\vec{q} \in \mathbb{Z}^{n}}\right| \vec{q}^{2} g_{\vec{q}}\right|^{2}<\infty\right\} ;
$$

(iii) the resolvent $\left(\mathrm{H}_{0}(\vec{k})-\mu\right)^{-1}$ of $\mathrm{H}_{0}(\vec{k})$ is a compact operator for all $\mu \notin \sigma\left(\mathrm{H}_{0}(\vec{k})\right)$, where $\sigma\left(\mathrm{H}_{0}(\vec{k})\right)$ is the spectrum of $\mathrm{H}_{0}(\vec{k})$.

Proof. - (i) and (ii) are obvious, since:

$$
\left(\mathrm{H}_{0} f\right)(\vec{\xi})=\vec{\xi}^{2} \hat{f}(\vec{\xi}) .
$$

(iii) The resolvent $\left(\mathrm{H}_{0}(\vec{k})-\mu\right)^{-1}$ is the multiplication operator by:

$$
\psi(\vec{q})=\left[(\vec{k}+\mathrm{E} \vec{q})^{2}-\mu\right]^{-1}
$$

Let $\chi_{M}$ be the characteristic function of the set $\left\{\vec{q} \in \mathbb{Z}^{n} \mid \vec{q}^{2} \leqq \mathrm{M}\right\}$ and $\mathrm{D}_{\mathrm{M}}$ the multiplication operator by $\psi(\vec{q}) \chi_{M}(\vec{q}) . \quad \mathrm{D}_{\mathrm{M}}$ is a compact (even nuclear) operator, and:

$$
\begin{equation*}
\left\|\left(\mathrm{H}_{0}(\vec{k})-\mu\right)^{-1}-\mathrm{D}_{\mathrm{M}}\right\|=\operatorname{Sup}_{\vec{q}>\mathrm{M}}\left[(\vec{k}+\mathrm{E} \vec{q})^{2}-\mu\right]^{-1} \rightarrow 0, \tag{10}
\end{equation*}
$$

as $M \rightarrow \infty$. Thus $\left(H_{0}(\vec{k})-\mu\right)^{-1}$ is compact as the uniform limit of the sequence $\left\{D_{M}\right\}$ of compact operators.

Let us denote by $\left\{\vec{v}_{\vec{q}}\right\}_{\vec{q} \in \mathbb{Z}^{n}}$ the Fourier coefficients of the periodic function $v$ :

$$
\begin{equation*}
\hat{v}_{\vec{q}}=\mathrm{L}^{-n / 2} \int_{\mathrm{C}^{n}} d x \exp (-i \mathrm{E} \cdot \vec{q} \cdot \vec{x}) v(\vec{x}) \tag{11}
\end{equation*}
$$

Notice that $v \in \mathrm{~L}^{p}\left(\mathrm{C}^{n}\right)$ for all $p \in[1, s]$. To establish the relation between the Fourier coefficients of $v$ and the operator $\mathrm{V}(\vec{k})$ we need the following result:

Lemma 3. - Given $\varphi, \psi: \mathbb{Z}^{n} \rightarrow \mathbb{C}$, we define an operator $\mathrm{A}_{\varphi \psi}: l_{n}^{2} \rightarrow l_{n}^{2}$ as follows:

$$
\left(\mathrm{A}_{\varphi \psi} g\right)_{\vec{q}}=\sum_{\vec{m} \in \mathbb{Z}^{n}} \varphi(\vec{m}) \psi(\vec{q}-\vec{m}) g_{\vec{q}-\vec{m}}
$$

Assume that $2 \leqq p<\infty, \psi \in l^{p}\left(\mathbb{Z}^{n}\right)$ and let $\{\varphi(\vec{q})\}$ be the Fourier coefficients of a function $\Phi$ belonging to $\mathrm{L}^{p}\left(\mathbb{C}^{n}\right)$. Then $\mathrm{A}_{\varphi \psi}$ is a compact operator and one has:

$$
\begin{equation*}
\left\|\mathrm{A}_{\varphi \psi}\right\| \leqq \mathrm{L}^{-(n / 2)-(n / p)}\|\Phi\|_{L^{p}\left(\mathrm{C}^{n}\right)}\|\psi\|_{l^{p}\left(\mathbb{Z}^{n}\right)} . \tag{12}
\end{equation*}
$$

Proof. - For $g=\left\{g_{\vec{q}}\right\} \in l_{n}^{2}$, define $\psi g=\left\{\psi(\vec{q}) g_{\vec{g}}\right\} . \quad$ By the Hölder inequality, $\psi g \in l_{n}^{r}$ with $r^{-1}=(1 / 2)+p^{-1}$, i.e. $1 \leqq r<2$, and:

$$
\|\psi g\|_{r} \leqq\|\psi\|_{p}\|g\|_{2}
$$

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Let:

$$
\gamma(x)=\mathrm{L}^{-n / 2} \sum_{\vec{q} \in \mathbb{Z}^{n}} \exp (i \mathrm{E} \vec{q} \cdot \vec{x}) \psi(\vec{q}) g_{\vec{q}}, \quad x \in \mathrm{C}^{n}
$$

By the Hausdorff-Young inequality [8], $\gamma \in \mathrm{L}^{r^{\prime}}\left(\mathrm{C}^{n}\right)$ with $\left(r^{\prime}\right)^{-1}=1-r^{-1}=1 / 2-p^{-1}$ and:

$$
\begin{equation*}
\|\gamma\|_{r^{\prime}} \leqq \mathrm{L}^{\left(n / r^{\prime}\right)-(n / 2)}\|\psi g\|_{r} \leqq \mathrm{~L}^{\left(n / r^{\prime}\right)-\left(n^{\prime} / 2\right)}\|\psi\|_{p}\|g\|_{2} \tag{13}
\end{equation*}
$$

Since $1 / 2=p^{-1}+\left(r^{\prime}\right)^{-1}$ and $\Phi \in \mathrm{L}^{p}\left(\mathrm{C}^{n}\right)$, the Hölder inequality implies that $\Phi \gamma \in \mathrm{L}^{2}\left(\mathrm{C}^{n}\right)$ and:

$$
\begin{equation*}
\|\Phi \gamma\|_{2} \leqq\|\Phi\|_{p}\|\gamma\|_{r^{\prime}} \leqq \mathrm{L}^{\left(n / r^{\prime}\right)-(n / 2)}\|\Phi\|_{p}\|\psi\|_{p}\|g\|_{2} \tag{14}
\end{equation*}
$$

Now:

$$
\left(\mathrm{A}_{\varphi \psi} g\right)_{\vec{q}}=\int_{\mathrm{C}^{n}} d x \exp (-i \mathrm{E} \cdot \vec{q} \cdot \vec{x}) \Phi(\vec{x}) \gamma(\vec{x})
$$

and by Plancherel's theorem we have:

$$
\begin{equation*}
\left\|\mathrm{A}_{\varphi \psi} g\right\|_{2}=\mathrm{L}^{n / 2}\|\Phi \gamma\|_{2} \leqq \mathrm{~L}^{n / r^{\prime}}\|\Phi\|_{p}\|\psi\|_{p}\|g\|_{2} \tag{15}
\end{equation*}
$$

This shows that $A_{\varphi \psi}$ is defined everywhere with the bound (12) :
(b) Let $\mathrm{D}_{\mathrm{M}}$ be the multiplication operator by $\psi_{\mathrm{M}}(\vec{q})=\psi(\vec{q}) \chi_{M}(\vec{q})$ (see the proof of Lemma 2). By $(a), \mathrm{A}_{\varphi \psi_{M}}$ is bounded, and $\mathrm{A}_{\varphi \psi_{M}}$ is non-zero only on a subspace of finite dimension. Therefore $A_{\varphi \psi_{M}}$ is nuclear. By using (12) we obtain:

$$
\begin{equation*}
\left\|\mathrm{A}_{\varphi \psi}-\mathrm{A}_{\varphi \Psi_{\mathrm{M}}}\right\| \leqq \mathrm{L}^{(n / 2)-(n / p)}\|\Phi\|_{p}\left\|\left(1-\chi_{\mathrm{M}}\right) \psi\right\|_{p} \tag{16}
\end{equation*}
$$

Since $\psi \in l_{n}^{p},\left\|\left(1-\chi_{\mathrm{M}}\right)\right\| \|_{p} \rightarrow 0$ as $\mathrm{M} \rightarrow \infty$. This proves the compactness of $\mathrm{A}_{\varphi \psi}$.
Lemma 4. - Let Y be the operator in $l_{n}^{2}$ defined by:

$$
\begin{equation*}
(\mathrm{Y} g)_{\vec{q}}=\mathrm{L}^{-n / 2} \sum_{\vec{m} \in \mathbb{Z}^{n}} \hat{v}_{\vec{m}} g_{\vec{q}-\vec{m}} \tag{17}
\end{equation*}
$$

Then:
(i) $\mathrm{D}_{0} \cong \mathrm{D}(\mathrm{Y})$ and Y is symmetric on $\mathrm{D}_{0}$;
(ii) Y is relatively compact with respect to $\mathrm{H}_{0}(\vec{k})$;
(iii) $\mathrm{V}(\vec{k})=\mathrm{Y}$ on $\mathrm{D}_{0}$, for all $\vec{k} \in \Gamma_{n}$ (in particular $\mathrm{V}(\vec{k})$ is independent of $\left.\vec{k}\right)$;
(iv) $\mathrm{H}(\vec{k})=\mathrm{H}_{0}(\vec{k})+\mathrm{Y}$ and $\mathrm{D}(\mathrm{H}(\vec{k}))=\mathrm{D}_{0}$.

Proof. - (i) If $g \in \mathrm{D}_{0}$, then $g=[\mathrm{H}(\overrightarrow{0})+1]^{-1}$ for some $h \in l_{n}^{2}$. (15) shows that $\|\mathrm{Y} g\|_{2}<\infty$, therefore $\mathrm{D}_{0} \subseteq \mathrm{D}(\mathrm{Y})$. By using $\bar{v}_{-\vec{q}}=v_{\vec{q}}$, one obtains easily that $(f, \mathrm{Y} g)=(\mathrm{Y} f, g)$ for $f, g \in \mathrm{D}_{0}$;
(ii) $\mathrm{Y}\left(\mathrm{H}_{0}(\vec{k})+1\right)^{-1}$ is of the form $\mathrm{A}_{\varphi \psi}$, with $\Phi(x)=\mathrm{L}^{-n / 2} v(\vec{x})$ and $\psi(\vec{q})=\left[(k+\mathrm{E} \vec{q})^{2}+1\right]^{-1}$. Notice that $\psi \in l_{n}^{p}$ for each $p>n / 2$. As $v \in \mathrm{~L}^{s}\left(\mathrm{C}^{n}\right)$ for $s=2$ if $n=2,3$ and $s>n / 2$ if $n \geqq 4$, Lemma 3 implies that $\mathrm{Y}\left(\mathrm{H}_{0}(\vec{k})+1\right)^{-1}$ is compact;

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(iii) this can be verified by calculating the Fourier transform of $\mathrm{V} f$;
(iv) by (i) and (ii), $\mathrm{H}_{0}(\vec{k})$ is self-adjoint. $\mathrm{H}(\vec{k})=\mathrm{H}_{0}(\vec{k})+\mathrm{Y}$ follows from (iii) and Lemmas 1 and 2.

## 4. Proof of Theorem 3

Let $f$ be an eigenvector of H , i.e. $\mathrm{H} f=\lambda f$ for some $\lambda \in \mathbb{R}$. By defining $v^{\prime}(x)=v(x)-\lambda$ and $\mathrm{H}^{\prime}=\mathrm{H}_{0}+\mathrm{V}^{\prime}$, we have $\mathrm{H}^{\prime} f=0$. Since $\mathrm{V}^{\prime}$ satisfies also the hypothesis (2), it is possible to assume without loss of generality that $\lambda=0$.
Let $\Gamma_{0}=\left\{\vec{k} \in \Gamma \mid(\mathscr{U} f)(\vec{k}) \neq 0\right.$ in $\left.l_{n}^{2}\right\}$. $\Gamma_{0}$ is measurable. Since $\mathrm{H}(\vec{k})(\mathscr{U} f)(\vec{k})=0, \mathrm{H}(\vec{k})$ must have the eigenvalue 0 for almost all the $\vec{k} \in \Gamma_{0}$. We will show that, for all $p \in\left(k_{1}, \ldots, k_{n-1}, 0\right) \in \mathbb{R}^{n-1}$ the set $\theta(\vec{p})$ of the points $k_{n} \in(0, \mathrm{E})$ such that $0 \in \sigma\left(\mathrm{H}\left(\vec{p}+k_{n} \mathrm{E}^{-1} \vec{e}_{n}\right)\right)$ is a set of measure zero. Thus the measure of $\Gamma_{0}$ is zero, i.e. $(\mathscr{U} f)(\vec{k})=0$ a.e., i.e. $f=0$. Therefore H cannot have any eigenvalues.
Fix $\vec{p}=\left(\vec{k}_{1}, \ldots, \vec{k}_{n-1}\right)$. To show that the measure of $\theta(\vec{p})$ is zero, we shall use the Fredholm theory of holomorphic families of operators of type (A), [7]. Let $\Omega$ be the following complex domain:

$$
\begin{equation*}
\Omega=\{\mathscr{X}+i r \mid \mathscr{X} \in(0,1), r \in \mathbb{R}\} . \tag{18}
\end{equation*}
$$

For $z \in \Omega$, we define $\mathrm{H}_{0}\left(\vec{p}, z \vec{e}_{n}\right)$ to be the multiplication operator in $l_{n}^{2}$ by $\left(\vec{p}+z \vec{e}_{n}+\mathrm{E} \vec{q}\right)^{2}$ and:

$$
\begin{equation*}
\mathrm{H}\left(\vec{p}, z \vec{e}_{n}\right)=\mathrm{H}_{0}\left(\vec{p}, z \vec{e}_{n}\right)+\mathrm{Y} . \tag{19}
\end{equation*}
$$

We shall see that:
(I) $\left\{\mathrm{H}\left(\vec{p}, z \vec{e}_{n}\right)\right\}$ is a holomorphic family of type (A) with respect to $z$. (See the terminology in [7]);
(II) the resolvent of $\mathrm{H}\left(\vec{p}, z \vec{e}_{n}\right)$ is compact;
(III) the resolvent set of $\mathrm{H}\left(\vec{p}, z \vec{e}_{n}\right)$ is not empty.

Under these conditions, Theorem VII.1.10 of [7] says that we have the following alternative:

- either $0 \in \sigma\left(\mathrm{H}\left(\vec{p}, z \vec{e}_{n}\right)\right)$ for each $z \in \Omega$;
- or every compact $\Omega_{0}$ in $\Omega$ contains only a finite number of points $z$ such that $0 \in \sigma\left(\mathrm{H}\left(\vec{p}, z \vec{e}_{n}\right)\right)$.
We shall show that:
(IV) 0 belongs to the resolvent set of $\mathrm{H}\left(\vec{p}, z \vec{e}_{n}\right)$ for $\operatorname{Im} z$ sufficiently large. Hence the first alternative is excluded, so that the measure of $\theta(\vec{p})$ is zero.
The remainder of the paper is devoted to the verification of the properties I to IV of $\mathrm{H}\left(\vec{p}, z \vec{e}_{n}\right)$. To simplify the notations we write $\mathrm{H}(\vec{p}, \vec{z})$ for $\mathrm{H}\left(\vec{p}, z \vec{e}_{n}\right)$.

Lemma 5. - (i) $\mathrm{H}_{0}(\vec{p}, z)$ is a self-adjoint holomorphic family of type (A) in $\Omega$ with domain $\mathrm{D}\left(\mathrm{H}_{0}(\vec{p}, z)\right)=\mathrm{D}_{0}$;

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(ii) $\forall z \in \Omega$, the resolvent of $\mathrm{H}_{0}(\vec{p}, z)$ is compact;
(iii) 0 belongs to the resolvent set $\rho\left(\mathrm{H}_{0}(\vec{p}, z)\right)$ of $\mathrm{H}_{0}(\vec{p}, z)$ for all $z$ with $\operatorname{Im} z \neq 0$.

Proof. - (i) Let $\mathrm{P}_{j}(j=1, \ldots n)$ be the following operator in $l_{n}^{2}$ :

$$
\begin{equation*}
\mathrm{P}_{j} g_{\vec{q}}=g_{j} g_{\vec{q}} . \tag{20}
\end{equation*}
$$

One has:

$$
\begin{equation*}
\mathrm{H}_{0}(\vec{p}, z)=\left(\vec{p}+\mathrm{E} \overrightarrow{\mathrm{P}}+z \vec{e}_{n}\right)^{2}=(\vec{p}+\mathrm{E} \overrightarrow{\mathrm{P}})^{2}+\mathrm{E}^{2} z^{2}+2 \mathrm{E}^{2} z \mathrm{P}_{n}, \tag{21}
\end{equation*}
$$

and the result is immediate:
(ii) the proof is the same as in Lemma 2 (iv).
(iii) for $z=\mathscr{X}+i r$, we have:

$$
\begin{equation*}
\operatorname{Im}\left(\vec{p}+\mathrm{E} \vec{q}+z \vec{e}_{n}\right)^{2}=2 \mathrm{E}^{2} r\left(\mathscr{X}+q_{n}\right) \tag{22}
\end{equation*}
$$

which is different from zero if $r \neq 0$. Since $q_{n} \in \mathbb{Z}$ and $\mathscr{X} \in(0,1)$ it follows that:

$$
\left\|\left[\mathrm{H}_{0}(\vec{p}, z)\right]^{-1}\right\|=\operatorname{Sup}_{\vec{q} \in \mathbb{Z}^{n}}\left|\left(\vec{p}+\mathrm{E} \vec{q}+z \vec{e}_{n}\right)^{2}\right|^{-1}<\infty
$$

i. e. $0 \in \rho\left(\mathrm{H}_{0}(\vec{p}, z)\right)$.

Lemma 6. - (i) $\mathrm{H}(\vec{p}, z)$ is a self-adjoint holomorphic family of type $(\mathrm{A})$ in $\Omega$ with domain $\mathrm{D}_{0}$;
(ii) $\forall z \in \Omega$ the resolvent of $\mathrm{H}(\vec{p}, z)$ is compact;
(iii) for all $\vec{p} \in \Gamma^{n-1}$ and $z \in \Omega, \rho(\mathrm{H}(\vec{p}, z))$ is not empty.

Proof. - (i) this follows from Lemmas 5 (i) and 4 (ii);
(iii) it suffices to show:

$$
\begin{equation*}
\lim _{\lambda \rightarrow+\infty}\left\|\mathrm{Y}\left[\mathrm{H}_{0}(\vec{p}, z)-i \lambda\right]^{-1}\right\|=0 \tag{23}
\end{equation*}
$$

since then the Neumann series for $[\mathrm{H}(\vec{p}, z)-i \lambda]^{-1}$, i.e.:

$$
\begin{equation*}
[\mathrm{H}(\vec{p}, z)-i \lambda]^{-1}=\left[\mathrm{H}_{0}(\vec{p}, z)-i \lambda\right]^{-1} \sum_{n=0}^{\infty}\left\{-\mathrm{Y}\left[\mathrm{H}_{0}(\vec{p}, z)-i \lambda\right]^{-1}\right\}^{n} \tag{24}
\end{equation*}
$$

is convergent if $\lambda$ is sufficiently large. Now, by (12):

$$
\begin{equation*}
\left\|\mathrm{Y}\left[\mathrm{H}_{0}(\vec{p}, z)-i \lambda\right]^{-1}\right\| \leqq \mathrm{L}^{-n / s}\|v\|_{s}\left\{\sum_{\vec{q} \in \mathbb{Z}^{n}}\left|\left(\vec{p}+\mathrm{E} \vec{q}+z \vec{e}_{n}\right)^{2}-i \lambda\right|^{-s}\right\}^{1 / s} \tag{25}
\end{equation*}
$$

We have with the notations $z=\mathscr{X}+\operatorname{ir}, \vec{k}=\left(\vec{p}, \mathscr{X} \vec{e}_{n}\right) \in \Gamma^{n}$ :

$$
\begin{align*}
\left|\left(\vec{p}+\mathrm{E} \vec{q}+z \vec{e}_{n}\right)^{2}-i \lambda\right|^{-2} & \leqq\left\{\left[(\vec{k}+\mathrm{E} \vec{q})^{2}-\mathrm{E}^{2} r^{2}\right]^{2}\right.  \tag{26}\\
& \left.+4 \mathrm{E}^{4} r^{2}\left[\mathscr{X}+q_{n}-\lambda\left(2 \mathrm{E}^{2} r\right)^{-1}\right]^{2}\right\}^{-1} \leqq\left[(\vec{k}+\mathrm{E} \vec{q})^{2}-\mathrm{E}^{2} r^{2}\right]^{-2} .
\end{align*}
$$

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This shows that each term of the sum in (26) converges to zero as $\lambda \rightarrow+\infty$, and that the series in (26) is uniformly majorized in $\lambda$ by a convergent serie (since $s>n / 2$ ). Therefore (23) is proven.
(If $z$ is such that $(\vec{k}+\mathrm{E} \vec{q})^{2}-\mathrm{E}^{2} r^{2}=0$ for certain $\vec{q} \in \mathbb{Z}^{n}$, then there exist $c>0$ and $\lambda_{0}<\infty$ such that $4 \mathrm{E}^{4} r^{2}\left[\mathscr{X}+q_{n}-\lambda\left(2 \mathrm{E}^{2} r\right)^{-1}\right]^{2} \geqq c$ for all these $\vec{q}$ and for each $\lambda \geqq \lambda_{0}$. For these values of $\vec{q}$ we can take as majorization in (26) the number $c^{-1}$ ).
(ii) Now we use the first and the second resolvent equation:

$$
\begin{gather*}
{[\mathrm{H}(\vec{p}, z)-\xi]^{-1}=[\mathrm{H}(\vec{p}, z)-\mu]^{-1}+(\xi-\mu)[\mathrm{H}(\vec{p}, z)-\xi]^{-1}[\mathrm{H}(\vec{p}, z)-\mu]^{-1}}  \tag{27}\\
{[\mathrm{H}(\vec{p}, z)-\mu]^{-1}=\left[\mathrm{H}_{0}(\vec{p}, z)-\mu\right]^{-1}-[\mathrm{H}(\vec{p}, z)-\mu]^{-1} \mathrm{Y}\left[\mathrm{H}_{0}(\vec{p}, z)-\mu\right]^{-1}}
\end{gather*}
$$

(27) shows that if $[H(\vec{p}, z)-\mu]^{-1}$ is compact for $\mu \in \rho(\mathrm{H}(\vec{p}, z))$ then $[\mathrm{H}(\vec{p}, z)-\xi]^{-1}$ is compact for each $\xi \in \rho(\mathrm{H}(\vec{p}, z))$. Since $\left[\mathrm{H}_{0}(\vec{p}, z)-\mu\right]^{-1}$ and $\mathrm{Y}\left[\mathrm{H}_{0}(\vec{p}, z)-\mu\right]^{-1}$ are compact if $\mu \in \rho\left(\mathrm{H}_{0}(\vec{p}, z)\right)$, by (28) it suffices to show that:

$$
\rho\left(\mathrm{H}_{0}(\vec{p}, z)\right) \cap \rho(\mathrm{H}(\vec{p}, z)) \neq \varnothing
$$

We know from (iii) that there exists a point $\mu_{0} \in \rho(\mathrm{H}(\vec{p}, z))$. If $\mu_{0} \notin \rho\left(\mathrm{H}_{0}(\vec{p}, z)\right)$, there exists a point close to $\mu \in \rho\left(\mathrm{H}_{0}(\vec{p}, z)\right) \cap \rho(\mathrm{H}(\vec{p}, z))$, since:
( $\alpha$ ) $\rho(\mathrm{H}(\vec{p}, z))$ is open;
$(\beta) \sigma\left(\mathrm{H}_{0}(\vec{p}, z)\right)$ consists of isolated eigenvalues only, because the resolvent of $\mathrm{H}_{0}(\vec{p}, z)$ is compact ([7]. Thm. III 6.29).

By Lemma 6 we have verified the properties (I) to (III) of the family $\{\mathrm{H}(\vec{p}, z)\}$. It now remains to prove (IV) i. e. $0 \in \rho(\mathrm{H}(\vec{p}, z))$ for some $z=\mathscr{X}+i r$ in $\Omega$. We have seen that $0 \in \rho\left(\mathrm{H}_{0}(\vec{p}, z)\right)$ if $r \neq 0$. We shall show that:

$$
\begin{equation*}
\lim _{r \rightarrow \infty}\left\|\mathrm{Y}\left[\mathrm{H}_{0}(\vec{p}, \mathscr{X}+i r)\right]^{-1}\right\|=0 \tag{29}
\end{equation*}
$$

By using the Neumann series (24) with $\lambda=0$ and $r$ sufficiently large, (29) implies $0 \in \rho(\mathrm{H}(\vec{p}, z))$ if $r=\operatorname{Im} z$ is sufficiently large.

To obtain (29), we use the inequality (25). By virtue of the first inequality in (26), it suffices to show that:

$$
\begin{equation*}
\lim _{r \rightarrow \infty} \sum_{\vec{q} \in \mathbb{Z}^{\prime \prime}}\left\{\left[\left(\vec{q}+\frac{\vec{k}}{\mathrm{E}}\right)^{2}-r^{2}\right]^{2}+4 r^{2}\left|q_{n}+\mathscr{X}\right|^{2}\right\}^{-s / 2}=0, \tag{30}
\end{equation*}
$$

which will be done in the next section.

## 5. Estimation of the series (30)

We now show that (30) holds if $s=2$ for $n=2,3, s>n-2$ for $n \geqq 4$ and $\mathscr{X} \in(0,1)$. We use the following notations:

$$
\begin{equation*}
a=2 r\left|q_{n}+\mathscr{X}\right|, \quad b=\left(q_{n}+\mathscr{X}\right)^{2}-r^{2} \tag{31}
\end{equation*}
$$

$$
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$$

We set $\vec{p}=\mathrm{E}^{-1}\left(k_{1}, \ldots, k_{n-1}\right) \in \Gamma_{1}^{n-1}$, where $\Gamma_{1}^{n-1}=\left\{\vec{p} \in \mathbb{R}^{n-1} \mid 0 \leqq p_{j}<1\right\}$, and:

$$
\begin{equation*}
\mathrm{S}\left(q_{n}, r\right)=\sum_{\vec{m} \in \mathbb{Z}^{n-1}}\left\{\left[(\vec{m}+\vec{p})^{2}+b\right]^{2}+a^{2}\right\}^{-s / 2} \tag{32}
\end{equation*}
$$

(30) is then equivalent to:

$$
\begin{equation*}
\operatorname{Lim}_{r \rightarrow \infty} \sum_{q_{n} \in \mathbb{Z}} \mathrm{~S}\left(q_{n}, r\right)=0 \tag{33}
\end{equation*}
$$

To prove (33), we first give a preliminary estimate in Lemma 7.
Lemma 7. - Let $\delta>0, c>0$ and $\mathrm{R}>0$. Then:

$$
\begin{equation*}
\varepsilon \equiv \inf _{r \geqq \mathrm{R}} \inf _{\substack{a \geqq \delta r \\ b \geqq-r^{2} \\ t, z \geqq \\|t-z| \leqq c}} \frac{\left(z^{2}+b\right)^{2}+a^{2}}{\left(t^{2}+b\right)^{2}+a^{2}}>0 \tag{34}
\end{equation*}
$$

Proof. - Setting $\alpha=a / r, \beta=b r^{-2}, \sigma=z / r, \tau=t / r$ and $\Omega_{r}=\{(\alpha, \beta, \sigma, \tau) \mid \alpha \geqq \delta, \beta \geqq-1$, $\left.\sigma \geqq 0, \tau \geqq 0,|\sigma-\tau| \leqq c r^{-1}\right\}$, we see that (34) is equivalent to:

$$
\begin{equation*}
\varepsilon=\inf _{r \geqq \mathrm{R} \Omega_{\Omega_{r}}} \frac{\left(\sigma^{2}+\beta\right)^{2}+(\alpha / r)^{2}}{\left(\tau^{2}+\beta\right)^{2}+(\alpha / r)^{2}}>0 \tag{35}
\end{equation*}
$$

The quotient on the r.h.s. of (35) is $\geqq 1$ if $\left|\tau^{2}+\beta\right| \leqq\left|\sigma^{2}+\beta\right|$. Hence the infimum is obtained by taking $\left|\tau^{2}+\beta\right| \geqq\left|\sigma^{2}+\beta\right|$. Under this restriction we have:

$$
\begin{equation*}
\frac{\left(\sigma^{2}+\beta\right)^{2}+(\alpha / r)^{2}}{\left(\tau^{2}+\beta\right)^{2}+(\alpha / r)^{2}} \geqq \max \left[\frac{\left(\sigma^{2}+\beta\right)^{2}}{\left(\tau^{2}+\beta\right)^{2}}, \frac{\left(\sigma^{2}+\beta\right)^{2}+(\alpha / r)^{2}}{\left(\tau^{2}+\beta\right)^{2}+2(\alpha / r)^{2}}\right] \tag{36}
\end{equation*}
$$

Also notice the following inequalities, valid on each $\Omega_{r}$ with $r \geqq \mathrm{R}$ :

$$
\begin{gather*}
\left.\tau^{2}+\beta=[(\tau-\sigma)+\sigma)\right]^{2}+\beta \leqq 2(\tau-\sigma)^{2}+2 \sigma^{2}+\beta  \tag{37}\\
=2\left(\sigma^{2}+\beta\right)-\beta+2(\tau-\sigma)^{2} \leqq 2\left(\sigma^{2}+\beta\right)+1+2 c^{2} \mathrm{R}^{-2} \\
\left|\left(\sigma^{2}+\beta\right)-\left(\tau^{2}+\beta\right)\right| \leqq(\sigma+\tau)|\sigma-\tau| \leqq(\sigma+\tau) c r^{-1} \tag{38}
\end{gather*}
$$

$$
\left(\tau^{2}+\beta\right)^{2} \leqq 2\left(\sigma^{2}+\beta\right)^{2}+2(\sigma+\tau)^{2} c^{2} r^{-2}
$$

We denote by $\varepsilon_{+}$and $\varepsilon_{-}$the infimum in (35) under the restriction $\sigma^{2}+\beta \geqq 1$ and $\sigma^{2}+\beta \in[-1,+1]$ respectively. It suffices to show that $\varepsilon_{+}>0$ and $\varepsilon_{-}>0$. In the first case (i.e. for $\sigma^{2}+\beta \geqq 1$ ), we use the first expression on the r.h.s. of (36) and the inequality (37). Setting $x=\sigma^{2}+\beta$, we see that:

$$
\begin{equation*}
\varepsilon_{+}=\inf _{x \geqq 1} \frac{x^{2}}{\left(2 x+1+2 c^{2} \mathrm{R}^{-2}\right)^{2}}>0 \tag{40}
\end{equation*}
$$

In the second case (i.e. for $\sigma^{2}+\beta \in[-1,+1]$ ), we have $\sigma^{2} \leqq 2$, hence $\sigma+\tau \leqq 2 \sqrt{2}+c \mathbf{R}^{-2} \equiv \eta$. After inserting this into (39) and using the second expression on the r.h.s. of (36), one obtains by setting $y=\left(\sigma^{2}+\beta\right)^{2}$ :

$$
\begin{align*}
& \varepsilon_{-}=\inf _{r \geqq R} \inf _{\substack{0 \leqq y \leqq 1 \\
\alpha \geqq \delta}} \frac{y+(\alpha / r)^{2}}{2 y+2 \eta^{2} c^{2} r^{-2}+2(\alpha / r)^{2}}  \tag{41}\\
&=\inf _{r \geqq R \inf _{\alpha \geqq i} 2 \eta^{2} c^{2} r^{-2}+2(\alpha / r)^{2}}=\frac{\delta^{2}}{2 \eta^{2} c^{2}+2 \delta^{2}}>0 .
\end{align*}
$$

Proof of (33). - Let $\vec{m} \in \mathbb{Z}^{n-1}$ and $\Gamma(\vec{m})$ be the cube:

$$
\Gamma(\vec{m})=\left\{\vec{x} \in \mathbb{R}^{n-1} \mid \vec{x}=\vec{p}+\vec{m}+\vec{y}, \vec{y} \in \Gamma_{1}^{n-1}\right\}
$$

We have $\Gamma(\vec{m}) \cap \Gamma\left(\vec{m}^{\prime}\right)=\emptyset$ if $\vec{m} \neq \vec{m}^{\prime}$ and:

$$
\mathbb{R}^{n-1}=\bigcup_{\vec{m} \in \mathbb{Z}^{n-1}} \Gamma(\vec{m})
$$

Let $c=\sqrt{n-1}$. Then for each $\vec{x} \in \Gamma(\vec{m})$ and each $\vec{m} \in \mathbb{Z}^{n-1}$ :

$$
||\vec{m}+\vec{p}|-|\vec{x}|| \leqq c
$$

Let $\delta=1 / 2 \min (\mathscr{X}, 1-\mathscr{X})$. By assumption $\delta>0$; since $a \geqq \dot{\delta} r$ and $b \geqq-r^{2}$, Lemma 7 implies the existence of a number $\varepsilon>0$ such that, for each $\vec{m} \in \mathbb{Z}^{n-1}$, each $x \in \Gamma(\vec{m})$, each $a \geqq \delta r$ and $b \geqq-r^{2}$ and all $r \geqq \mathrm{R}$ :

$$
\begin{equation*}
\left[(\vec{m}+\vec{p})^{2}+b^{2}\right]+a^{2} \geqq \varepsilon\left[\left(\vec{x}^{2}+b\right)^{2}+a^{2}\right] . \tag{42}
\end{equation*}
$$

Thus:

$$
\begin{align*}
\mathrm{S}\left(q_{n}, r\right)= & \sum_{\vec{m} \in \mathbb{Z}^{n-1}}\left\{\left[(\vec{m}+\vec{p})^{2}+b^{2}\right]+a^{2}\right]^{-s / 2}  \tag{43}\\
= & \sum_{\vec{m} \in \mathbb{Z}^{n-1}} \int_{\Gamma(\vec{m})} d x\left\{\left[(\vec{m}+\vec{p})^{2}+b\right]^{2}+a^{2}\right]^{-s / 2} \\
& \leqq \varepsilon^{-1} \sum_{\vec{m}} \int_{\Gamma(\vec{m})} d x\left\{\left(\vec{x}^{2}+b\right)^{2}+a^{2}\right\}^{-s / 2} \\
& =\varepsilon^{-1} \int_{\mathbb{R}^{n-1}} d x\left\{\left(\vec{x}^{2}+b\right)^{2}+a^{2}\right\}^{-s / 2} \\
& =\frac{1}{2} \varepsilon^{-1} w_{n-1} \int_{0}^{\infty} y^{(n-3) / 2}\left\{(y+b)^{2}+a^{2}\right\}^{-s / 2} d y
\end{align*}
$$

where we have introduced spherical polar coordinates, $y=|\vec{x}|^{2}$ and $w_{n-1}$ denotes the area of the unit sphere in $\mathbb{R}^{n-1}$.

$$
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$$

To estimate the integral in (43), we distinguish the two cases $b \geqq 0$ and $b<0 . \quad$ For $b \geqq 0$, we have $\left\{(y+b)^{2}+a^{2}\right\}^{-s / 2} \leqq\left\{y^{2}+a^{2}+b^{2}\right\}^{-s / 2}$, and (43) leads to:

$$
\mathrm{S}\left(q_{n}, r\right) \leqq \frac{1}{2} \varepsilon^{-1} w_{n-1}\left(a^{2}+b^{2}\right)^{-s / 2+(n-1) / 4} \int_{0}^{\infty} z^{(n-3) / 2}\left(z^{2}+1\right)^{-s / 2} d z
$$

Notice that the integral in this expression is convergent since $s>n / 2$. By observing that:

$$
\begin{equation*}
a^{2}+b^{2}=\left[(q+\mathscr{X})^{2}+r^{2}\right]^{2} . \tag{44}
\end{equation*}
$$

we obtain:

$$
\begin{equation*}
\sum_{\left|q_{n}+\mathscr{X}\right| \geqq r} \mathrm{~S}\left(q_{n}, r\right) \leqq \mathrm{Cte} \sum_{\left|q_{n}+\mathscr{X}\right| \geqq r}\left|q_{n}+\mathscr{X}\right|^{-2 s+n-1} \tag{45}
\end{equation*}
$$

The hypothesis $s>n / 2$ implies that the last series is convergent so that this term tends to zero as $r \rightarrow \infty$.

We now turn to the case $b<0$. We set $z=(y+b) / a$. (43) then gives:

$$
\begin{equation*}
\mathrm{S}\left(q_{n}, r\right) \leqq \frac{1}{2} \varepsilon^{-1} w_{n-1} a^{-s+1} \int_{b / a}^{+\infty}(a z-b)^{(n-3) / 2}\left\{1+z^{2}\right\}^{-s / 2} d z \tag{46}
\end{equation*}
$$

If $n \geqq 3$, this leads to:

$$
\begin{align*}
& \mathrm{S}\left(q_{n}, r\right) \leqq c_{1} a^{-s+1} \int_{-\gamma}^{+\infty}\left[|a z|^{(n-3) / 2}+|b|^{(n-3) / 2}\right]\left\{1+z^{2}\right\}^{-s / 2} d z  \tag{47}\\
& \leqq c_{2} a^{-s+1}\left[|a|^{(n-3) / 2}+|b|^{(n-3) / 2}\right] \leqq c_{3} a^{-s+1}\left(a^{2}+b^{2}\right)^{(n-3) / 4} .
\end{align*}
$$

Using (47), (44) and (31), we obtain in this case that:

$$
\sum_{\left|q_{n}+X\right|<r} \mathrm{~S}\left(q_{n}, r\right) \leqq c_{4} r^{-s+1} r^{n-3} \sum_{\left|q_{n}+X\right|<r}\left|q_{n}+\mathscr{X}\right|^{-s+1}=\mathcal{O}\left(r^{-s+n-2} \log r\right)
$$

since $s \geqq 2$. Under the hypothesis $s>n-2$, this converges to zero as $r \rightarrow \infty$.
Finally, if $n=2$, one may bound the integral in (46) by a constant which is independent of $a$ and $b$ on the set $\left\{a \geqq a_{0}>0, b<0\right\}$; this is easily achieved by splitting the domain of integration into $\{z \mid a z-b \leqq 1\} \cup\{z \mid a z-b>1\}$. Thus:

$$
\mathrm{S}\left(q_{n}, r\right) \leqq c_{5} r^{-s+1}\left|q_{n}+\mathscr{X}\right|^{-s+1}, \quad \forall q_{n}, \quad \forall r \geqq r_{0}
$$

For any $s>3 / 2$, this implies that:

$$
\lim _{r \rightarrow \infty} \sum_{\left|q_{n}+X\right|<r} \mathrm{~S}\left(q_{n}, r\right)=0
$$

Remark 3. - One sees from the preceding proof that, for $n=3$, the limit in (33) is zero under the weaker hypothesis that $s>3 / 2$. By using a modified resolvent equation, one obtains the result of Theorem 1 for $s>3 / 2$. The case $s=2, n=3$ was first treated by Thomas
in [12]. Similarly, for $n=2$, a more careful estimate of the integral in (46) shows that it suffices to require $s>1$.

Remark 4. - Theorem 3 remains true if the condition of ortho-periodicity of $v$ is replaced by the weaker condition of periodicity. Indeed, the estimation of the series given in section 5 , may be applied if, instead of $\vec{a}_{i} \cdot \vec{a}_{j}=\delta_{i j}$, one requires only that $\vec{a}_{i} \cdot \vec{a}_{n}=\delta_{i n}$ (i.e. the vector $\vec{a}_{n}$ is orthogonal to the hyperplane spanned by $\left.\vec{a}_{1}, \ldots, \vec{a}_{n-1}\right)$. Clearly the direction $\vec{a}_{n}$ is distinguished in our estimation. A similar result for an arbitrary periodic lattice is given in Theorem XIII. 100 of [11], under a more restrictive assumption on the local behaviour of the function $v(\vec{x})$ than that of Theorem 3.

Remark 5. - We also have the following result which generalizes Theorem 1:
Theorem $1^{\prime} .-$ Let $v \in \mathrm{~L}_{\mathrm{loc}}^{s}\left(\mathbb{R}^{n} \backslash \mathrm{~N}\right)$, where s satisfies $s=2$ if $n=1,2,3$ and $s>n-2$ if $n \geqq 4$, and where N is a closed set of measure zero. Let H be a self-adjoint extension of $\hat{\mathrm{H}}$, $\mathrm{D}(\hat{\mathrm{H}})=\mathrm{C}_{0}^{\infty}\left(\mathbb{R}^{n} \backslash \mathrm{~N}\right)$. Suppose that $f \in \mathrm{~L}^{2}\left(\mathbb{R}^{n}\right)$ satisfies $\mathrm{H} f=\lambda$ ffor some $\lambda \in \mathbb{P}$ and $\mathrm{E}(\mathrm{A}) f=f$ for some compact subset of $\mathbb{R}^{n} \backslash \mathrm{~N}$ (i.e.f is an eigenvector of H having compact support in $\left.\mathbb{R}^{n} \backslash \mathrm{~N}\right)$. Then $f=0$.

Proof. - One has $\chi_{\mathrm{A}}() v.(.) \in \mathrm{L}^{s}\left(\mathbb{R}^{n}\right) . \quad$ Let C be a cube in $\mathbb{R}^{n}$ such that $\mathrm{A} \subseteq \mathrm{C}$. Define $w$ by:

$$
w\left(\vec{x}+\sum q_{i} \vec{a}_{i}\right)=\chi_{\mathrm{A}}(\vec{x}) v(\vec{x}), \quad \vec{x} \in \mathrm{C}
$$

$w$ is ortho-periodic and in $\mathrm{L}_{\text {loc }}^{s}\left(\mathbb{R}^{n}\right) . \quad$ Since $\left(\mathrm{H}_{0}+w\right) f=\lambda f$, one has $f=0$ by Theorem 3.
Remark 6. - The hypothesis " $\Sigma$ bounded" in Theorem $1(b)$ is essential. Assume for example that $v$ is such that $\mathrm{H}_{0}+v$ has pure point spectrum (e.g. $v(\vec{x}) \rightarrow+\infty$ as $|\vec{x}| \rightarrow \infty)$. Take $\Sigma=\mathbb{R}$. Then:

$$
\mathrm{F}(\Sigma) \mathscr{H}=\mathscr{H} \quad \text { and } \quad \mathrm{E}(\mathrm{~A}) \mathscr{H} \cap \mathrm{F}(\Sigma) \mathscr{H} \cap \mathscr{H}_{p}(\mathrm{H})=\mathrm{E}(\mathrm{~A}) \mathscr{H} .
$$

Since $\mathrm{E}(\mathrm{A}) \mathscr{H} \neq\{0\}$ if A has positive measure, it is clear that one cannot have $\mathrm{E}(\mathrm{A}) \mathscr{H} \cap \mathrm{F}(\Sigma) \mathscr{H} \cap \mathscr{H}_{p}(\mathrm{H})=\{0\}$ in this case.

Remark 7. - By combining our Theorem 1 with Proposition 4 of [2], one may also prove that $\mathrm{E}(\mathrm{A}) \mathscr{H} \cap \mathrm{F}(\Sigma) \mathscr{H}=\{0\}$ under assumptions of Theorem $1(b)$.

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