# The Bourgain Spaces and Recovery of Magnetic and Electric Potentials of Schrödinger Operators 

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## Recommended Citation

Zhang, Yaowei, "The Bourgain Spaces and Recovery of Magnetic and Electric Potentials of Schrödinger Operators" (2016). Theses and Dissertations--Mathematics. 41.
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Yaowei Zhang, Student<br>Dr. Russell Brown, Major Professor<br>Dr. Peter Hislop, Director of Graduate Studies

The Bourgain spaces and recovery of magnetic and electric potentials of Schrödinger operators

## DISSERTATION

A dissertation submitted in partial fulfillment of the requirements for the degree of Doctor of Philosophy in the College of Arts and Sciences at the University of Kentucky

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Director: Dr. Russell Brown, Professor of Mathematics
Lexington, Kentucky 2016

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## ABSTRACT OF DISSERTATION

The Bourgain spaces and recovery of magnetic and electric potentials of Schrödinger operators

We consider the inverse problem for the magnetic Schrödinger operator with the assumption that the magnetic potential is in $C^{\lambda}$ and the electric potential is of the form $p_{1}+\operatorname{div} p_{2}$ with $p_{1}, p_{2} \in C^{\tilde{\lambda}}$. We use semiclassical pseudodifferential operators on semiclassical Sobolev spaces and Bourgain type spaces. The Bourgain type spaces are defined using the symbol of the operator $h^{2} \Delta+h \mu \cdot D$. Our main result gives a procedure for recovering the curl of the magnetic field and the electric potential from the Dirichlet to Neumann map. Our results are in dimension three and higher.

KEYWORDS: Schrödinger operator, Pseudodifferential operators, Bourgain type spaces

The Bourgain spaces and recovery of magnetic and electric potentials of Schrödinger operators

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Date:
August 2, 2016

## ACKNOWLEDGMENTS

For the past years in mathematics department of UK, I have received support and encouragement from a lot of professors and students. Dr. Russell Brown is an advisor and friend. He can always think out a way to solve problems or find the helpful articles. I would like to thank the other members of my dissertation committee: James Mcdonough, Peter Perry, Mihai Tohoneanu, and Thomas Troland for their help.

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## Chapter 1 Introduction

Let $\Omega \subset \mathbf{R}^{n}, n \geq 3$ be a bounded domain with $C^{1,1}$ boundary. The magnetic Schrödinger operator is

$$
\begin{equation*}
H_{W, p}=\sum_{j=1}^{n}\left(D_{j}+W_{j}\right)^{2}+p \tag{1.1}
\end{equation*}
$$

where $D_{j}=\frac{1}{i} \frac{\partial}{\partial x_{j}}, W$ is the magnetic potential, and $p$ is the electric potential. We assume that $W \in C_{c}^{\lambda}\left(\mathbf{R}^{n} ; \mathbf{C}^{n}\right), p=p_{1}+\operatorname{div} p_{2}$ with $p_{1}, p_{2} \in C_{c}^{\tilde{\lambda}}\left(\mathbf{R}^{n} ; \mathbf{C}\right)$. We are using $C_{c}^{\alpha}\left(\mathbf{R}^{n}\right)$ to denote the space of functions which are Hölder continuous of exponent $\alpha$ and compactly supported in $\mathbf{R}^{n}$.

If we assume 0 is not a Dirichlet eigenvalue of $H_{W, p}$, then the boundary value problem

$$
\begin{cases}H_{W, p} u=0 & \text { in } \Omega \\ u=f & \text { on } \partial \Omega\end{cases}
$$

has a unique solution $u=u_{f} \in H^{1}(\Omega)$ for any $f \in H^{1 / 2}(\partial \Omega)$.
We will use the Dirichlet to Neumann map (DN map) to describe our boundary measurements. The DN map is defined by

$$
\Lambda_{W, p}:\left.f \rightarrow \frac{\partial u_{f}}{\partial \nu}\right|_{\partial \Omega}+i(W \cdot \nu) f
$$

where $\nu$ is the outer unit normal to $\partial \Omega$. Furthermore, for $f, g \in H^{1 / 2}(\partial \Omega)$, if we assume $H_{W, p} u_{f}=0$ in $\Omega, u_{f}=f$ on $\partial \Omega, H_{-W, p} v_{g}=0$ in $\Omega, v_{g}=g$ on $\partial \Omega, \varphi_{f}, \varphi_{g}$ are any functions in $H^{1}(\Omega)$ with $\varphi_{f}=f, \varphi_{g}=g$ on $\partial \Omega$, then the weak formulation for $\Lambda_{W, p}$ is

$$
\begin{equation*}
\left\langle\Lambda_{w, p} f, g\right\rangle=\int_{\Omega}\left(\nabla u_{f} \cdot \nabla \varphi_{g}+i W \cdot\left(u_{f} \nabla \varphi_{g}-\varphi_{g} \nabla u_{f}\right)+(W \cdot W) u_{f} \varphi_{g}\right) d x+\left\langle p u_{f}, \varphi_{g}\right\rangle \tag{1.2}
\end{equation*}
$$

Using the adjoint of $\Lambda_{W, p}$, we also have

$$
\begin{equation*}
\left\langle\Lambda_{w, p} f, g\right\rangle=\int_{\Omega}\left(\nabla \varphi_{f} \cdot \nabla v_{g}+i W \cdot\left(\varphi_{f} \nabla v_{g}-v_{g} \nabla \varphi_{f}\right)+(W \cdot W) \varphi_{f}, v_{g}\right) d x+\left\langle p \varphi_{f}, v_{g}\right\rangle \tag{1.3}
\end{equation*}
$$

From this definition it follows that $\Lambda_{W, p}$ is bounded on $H^{1 / 2}(\partial \Omega) \rightarrow H^{-1 / 2}(\partial \Omega)$. Since $p$ is not a function, following Brown [4], we define the "multiplication by $p$ " by $\langle p u, v\rangle=\int p_{1} u v-\int p_{2} \cdot \nabla(u v)$ when $u$ and $v$ are sufficently smooth.

The inverse problem for the magnetic Schrödinger operator is the problem of recovering curl $W$ and $p$ from $\Lambda_{W, p}$. This problem related to the inverse conductivity problem of Calderón [5]. Previous results for this problem concern different assumptions of $W, p$ (starting with Sylvester and Uhlmann[22] and later work in
[26],[27]). The inverse problem of the closely related magnetic Schrödinger equation, was first worked by Sun [21]. Nakamura, Sun and Uhlmann [15] considered $W$, $p \in C^{\infty}$ (Tolmasky [24] improved to $C^{1}$ and Salo [17] to Dini continuous). Salo [18] solved this in the case when $W$ is continuous with $D \cdot W \in L^{\infty}, p \in L^{\infty}$. Krupchyk and Uhlmann 12 worked on this problem when $W \in L^{\infty}$ and $q \in L^{\infty}$. Pohjola [16] studied the case when $W \in C^{2 / 3+\varepsilon}, p=p_{1}+\operatorname{div} p_{2}$ with $p_{1}, p_{2} \in C^{2 / 3+\varepsilon}$ for any small $\varepsilon$. Haberman [9] proves the case that on a ball $B \subset \mathbf{R}^{3}$ when $W$ is small in $W^{s, 3}$ for some $s>0$, and $p \in W^{-1,3}$.

We recover $p$ by using the complex geometrical optics (CGO) solutions of $H_{W, p} u=$ 0 in Bourgain type spaces. CGO solutions are solutions of the form

$$
u=e^{i x \cdot \zeta} v
$$

where $\zeta \in \mathbf{C}^{n}$, and satisfies $\zeta \cdot \zeta=0$. The use of CGO solutions for inverse problems first appeared in Sylvester-Uhlmann [22]. The Bourgain type spaces we use were introduced to the study of inverse problems by Haberman and Tataru [10]. The original Bourgain spaces were defined by Bourgain in [2]. The definition is as follows. In this definition and throughout this paper, we use $a \cdot b=\sum_{i} a_{i} b_{i}$ for vectors $a$ and $b$ in $\mathbf{C}^{n}$.

Definition 1.1. Let $q(\xi)=\xi \cdot \xi+2 \mu \cdot \xi$, where $\mu \in \mathbf{C}^{n}$ with $|\mu|=\sqrt{2}$ and $\mu \cdot \mu=0$, we define spaces $\dot{X}_{\mu, h}^{b}, X_{\mu, h}^{b}$ and $X_{\mu, h, \sigma}^{b}$ with the following norms

$$
\begin{aligned}
\|f\|_{\dot{X}_{\mu, h}^{b}} & =\left\||q(h \cdot)|^{b} \hat{f}(\cdot)\right\|_{2} \\
\|f\|_{X_{\mu, h}^{b}}^{b} & =\left\|(h+|q(h \cdot)|)^{b} \hat{f}(\cdot)\right\|_{2} \\
\|f\|_{X_{\mu, h, \sigma}^{b}}^{b} & =\left\|\left(h^{2(1-\sigma)}+|q(h \cdot)|^{2}\right)^{b / 2} \hat{f}(\cdot)\right\|_{2}
\end{aligned}
$$

for $f \in \mathcal{S}^{\prime}, \hat{f} \in L_{l o c}^{1}$ is a function, $b, h, \sigma \in \mathbf{R}$ with $|b|<1, \sigma \in[0,1), h>0$.
In this definition and below, we use $\mathcal{S}^{\prime}\left(\mathbf{R}^{n}\right)$ to denote the space of tempered distributions and $\mathcal{S}\left(\mathbf{R}^{n}\right)$ will denote the space of Schwartz functions.

Lemma 1.2. For $b<1$, the the spaces $\dot{X}_{\mu, h}^{b}, X_{\mu, h}^{b}$ and $X_{\mu, h, \sigma}^{b}$ are Banach spaces.
Proof. It is clear that the spaces $X_{\mu, h}^{b}$ and $X_{\mu, h, \sigma}^{b}$ are Banach spaces. The interesting point is to show that $\dot{X}_{\mu, h}^{b}$ is complete.

Suppose we have a Cauchy sequence $\left\{f_{n}\right\}$ in the space $\dot{X}_{\mu, h}^{b}$. Then by the definition of the space $\dot{X}_{\mu, h}^{b}, \hat{f}_{n}$ is an element of a weighted $L^{2}$-space where the norm of a function $g$ is $\|q(h \cdot) g\|_{L^{2}}$. Since the weighted $L^{2}$-space is complete, we can define

$$
g:=\lim _{n \rightarrow \infty} \hat{f}_{n} .
$$

We need to show that $g$ is a tempered distribution and thus that $g=\hat{f}$. Then we will have that $f=\lim _{n \rightarrow \infty} f_{n}$ where the limit occurs in the space $\dot{X}_{\mu, h}^{b}$. To show that $g$ is a tempered distribution, observe that we have

$$
\begin{equation*}
\int_{B(0, R)}|g(\xi)| d \xi \leq\left(\int_{B(0, R)}|g(\xi)|^{2}|q(h \xi)|^{2 b} d \xi\right)^{1 / 2}\left(\int_{B(0, R)}|q(h \xi)|^{-2 b} d \xi\right)^{1 / 2} \tag{1.4}
\end{equation*}
$$

From estimate (2.1), since $b<1$

$$
\int_{B(0, R)}|q(h \xi)|^{-2 b} d \xi \leq C R^{n-2 b} / h^{2 b}
$$

Using this in (1.4) implies that $g$ is a tempered distribution. This proves the Lemma.

To recover $p$, we need to prove the existence of CGO solutions in the space $\dot{X}_{\mu, h}^{1 / 2}$ space. We will need to define semiclassical pseudodifferential operators and define their action on Bourgain spaces in order to establish the existence of the CGO solutions. We need a variant of the pseudodifferential cutoff technique used by Takeuchi [23] and Kenig-Ponce-Vega [11]. Zworski [28] provides a good introduction to semiclassical pseudodifferential operators. To recover $W$, we follow the method of Salo [18, Section 6]. Even though our $W$ is less regular than in Salo, we are still able to adapt his proof. Once the CGO solutions are constructed, we can use a non-physical scattering transform to recover $p$. The main result is

Theorem 1.3. Let $\Omega \subset \mathbf{R}^{n}$, $n \geq 3$ be a bounded $C^{1,1}$ domain. If $W \in C^{\lambda}\left(\bar{\Omega} ; \mathbf{C}^{n}\right)$, $p=p_{1}+\operatorname{div}_{2}$ with $p_{1}, p_{2} \in C^{\tilde{\lambda}}(\bar{\Omega} ; \mathbf{C}) \cap H^{\lambda}\left(\mathbf{C}^{n}\right)$ and 0 is not a Dirichlet eigenvalue of $H_{W, p}$. Then when $\lambda \in(1 / 2,1), \tilde{\lambda} \in(0,1)$, curl $W$ is determined by $\Lambda_{W, p}$. Further, when $(\lambda+1) \tilde{\lambda}>\frac{3}{2}, 0<\tilde{\lambda}, \lambda<1, \Lambda_{W, p}$ determines $p$ uniquely.

The main interest of this result is the low regularity on $W$ and $p$. If we let $\lambda=\tilde{\lambda}$, then our theorem applies for $\lambda>(-1+\sqrt{7}) / 2 \approx 0.82$. Thus our theorem allows potentials which are not functions and represents an improvement over the work of Krupchyk and Uhlmann. The work of Pohjola assumes that $\lambda=\tilde{\lambda}$ and thus, our work is slightly more general. Unlike Haberman's recent work, there is not a restriction that $W$ is small. However, the result above is far from sharp.

Here is the structure of this dissertation. Chapter 2 talks about the properties of Bourgain type spaces and semiclassical pseudodifferential operators. In chapter 3 and chapter 4, we obtain the existence of CGO solutions in the space $\dot{X}_{\mu, h}^{1 / 2}$ and in semiclassical Sobolev spaces. We show the boundary value of CGO solutions can be determined by $\Lambda_{W, p}$ in chapter 5 . chapter 6 contains the recovery of curl $W$. We recover $p$ in chapter 7 . Chapter 8 talks about the future work.

## Chapter 2 Bourgain type spaces and semiclassical pseudodifferential operators

## Bourgain type spaces

We first present some properties of Bourgain type spaces, which are critical when we prove the existence of CGO solutions in the space $\dot{X}_{\mu, h}^{1 / 2}$.

Let $\langle x\rangle=\left(1+|x|^{2}\right)^{1 / 2}$, before we come to the proposition, we need this lemma.
Lemma 2.1. let $B_{r}(\eta)$ be the ball of radius $r$ with center $\eta, k \in(0,2), N$ is a large positive number, then

$$
\begin{equation*}
\int_{B_{r}(\eta)}\left|\frac{1}{q(h \xi)}\right|^{k} \leq C r^{n-k} / h^{k} \tag{2.1}
\end{equation*}
$$

and

$$
\begin{equation*}
\int\langle\eta-\xi\rangle^{-N}\left|\frac{h}{q(h \xi)}\right| \leq C \tag{2.2}
\end{equation*}
$$

for some constant $C$ depending only on the dimension $n$.
Proof. The proof of the first estimate (2.1) follows from R. Brown [3, page 82]. The second estimate is from the first estimate and is also in Haberman and Tataru [10, Lemma 2.2]. I will show how the first estimate works.

By a rotation in the variable $\xi$, it suffices to consider the case where $\mu=e_{1}+i e_{2}$. We define the zero set of $q(h \xi)$ by $\Sigma_{\mu}=\left\{\xi: q(h \xi)=h^{2}|\xi|^{2}+2 \mu \cdot(h \xi)=0\right\}$. We consider 3 cases of the ball $B_{r}(\eta)$.

Case 1: When $r<\frac{1}{100} h^{-1}$ and $\operatorname{dist}\left(\eta, \Sigma_{\mu}\right)<2 r$. We can rotate the variables $\left(\xi_{2}, \xi_{3}, \cdots, \xi_{n}\right)$ about the center of $\Sigma_{\mu}$, so that $B_{r}(\eta) \subset B_{3 r}(0)$. Since $q(h \xi)=|h \xi|^{2}+$ $2(\operatorname{Re} \mu+i \operatorname{Im} \mu) \cdot(h \xi)=|h \xi+\operatorname{Re} \mu|^{2}-|\operatorname{Re} \mu|^{2}+2 i \operatorname{Im} \mu \cdot(h \xi)$, we can define new variables $x_{1}=|h \xi+R e \mu|^{2}-|R e \mu|^{2}, x_{2}=\operatorname{Im} \mu \cdot(h \xi), x_{j}=h \xi_{j}, j=3,4, \cdots, n$. Then we have $\frac{d x_{1}}{d \xi_{j}}=2 h^{2} \xi_{j}+2 h(R e \mu)_{j}, \frac{d x_{2}}{d \xi_{j}}=2 h(\operatorname{Im} \mu)_{j}$. So we obtain that

$$
\begin{align*}
\int_{B_{r}(\eta)}\left|\frac{1}{q(h \xi)}\right|^{k} & \leq C h^{-n} \int_{B_{C h r}(0)}\left|\frac{1}{x_{1}+i x_{2}}\right|^{k} d x_{1} d x_{2} \cdots d x_{n} \\
& \leq C h^{-n}(h r)^{n-k}  \tag{2.3}\\
& \leq C r^{n-k} / h^{k} .
\end{align*}
$$

Case 2: When $\operatorname{dist}\left(\eta, \Sigma_{\mu}\right)>2 r$. Since $|q(h \xi)|$ is comparable to $h \operatorname{dist}\left(\xi, \Sigma_{\mu}\right)$ when $|\xi|<8 h^{-1}$ and comparable to $|h \xi|^{2}$ when $|\xi|>8 h^{-1}$. Thus

$$
\begin{equation*}
\frac{1}{|q(h \xi)|} \leq C \frac{1}{h r} . \tag{2.4}
\end{equation*}
$$

This gives

$$
\begin{equation*}
\int_{B_{r}(\eta)}\left|\frac{1}{q(h \xi)}\right|^{k} \leq C\left(\frac{1}{h r}\right)^{k} \int_{B_{r}(\eta)} 1 \leq C r^{n-k} / h^{k} \tag{2.5}
\end{equation*}
$$

Case 3: When $r>\frac{1}{100} h^{-1}$ and $\operatorname{dist}\left(\eta, \Sigma_{\mu}\right)<2 r$. We write $B_{r}(\eta)=B_{0} \cup B_{\infty}$, where $B_{0}=B_{r}(\eta) \cap B_{4 h^{-1}}(0)$ and $B_{\infty}=B_{r}(\eta) \backslash B_{4 h^{-1}}(0)$. By the argument in Case 1 and Case 2, we know

$$
\begin{equation*}
\int_{B_{0}}\left|\frac{1}{q(h \xi)}\right|^{k} \leq C / h^{n} \tag{2.6}
\end{equation*}
$$

We know that $\Sigma_{\mu} \subset B_{4 h^{-1}}(0)$, so

$$
\begin{equation*}
\frac{1}{|q(h \xi)|} \leq C \frac{1}{|h \xi|^{2}} \tag{2.7}
\end{equation*}
$$

on $B_{\infty}$. Thus

$$
\begin{equation*}
\int_{B_{\infty}}\left|\frac{1}{q(h \xi)}\right|^{k} \leq C \frac{1}{h^{2 k}} r^{n-2 k} \tag{2.8}
\end{equation*}
$$

which implies 2.1 since $r>\frac{1}{100} h^{-1}$.
Proposition 2.2. For any $\delta>0$, we have a constant $C=C(n, \delta)$ so that

$$
\begin{gather*}
\left\|\langle x\rangle^{-1 / 2-\delta} f\right\|_{X_{\mu, h}^{1 / 2}} \leq C\|f\|_{\dot{X}_{\mu, h}^{1 / 2}}  \tag{2.9}\\
\left\|\langle x\rangle^{-1 / 2-\delta} f\right\|_{\dot{X}_{\mu, h}^{-1 / 2}} \leq C\|f\|_{X_{\mu, h}^{-1 / 2}} \leq C h^{-1 / 2}\|f\|_{L^{2}}  \tag{2.10}\\
\left\|\langle x\rangle^{-1 / 2-\delta} f\right\|_{L^{2}} \leq C h^{-1 / 2}\|f\|_{\dot{X}_{\mu, h}^{1 / 2}}  \tag{2.11}\\
\|f\|_{\dot{X}_{\mu, h}^{1 / 2}} \leq\|f\|_{X_{\mu, h}^{1 / 2}} \leq\|f\|_{X_{\mu, h, \sigma}^{1 / 2}}  \tag{2.12}\\
\|f\|_{X_{\mu, h, \sigma}^{-1 / 2}} \leq\|f\|_{X_{\mu, h}^{-1 / 2}} \leq\|f\|_{\dot{X}_{\mu, h}^{-1 / 2}}  \tag{2.13}\\
\|f\|_{X_{\mu, h, \sigma}^{1 / 2}} \leq h^{-\sigma / 2}\|f\|_{X_{\mu, h}^{1 / 2}}  \tag{2.14}\\
\|f\|_{X_{\mu, h}^{-1 / 2}} \leq h^{-\sigma / 2}\|f\|_{X_{\mu, h, \sigma}^{-1 / 2}} \tag{2.15}
\end{gather*}
$$

The following proof uses ideas from Haberman and Tataru [10].
Proof. Let $\varphi(x)=\langle x\rangle^{-1 / 2-\delta}$, then

$$
\begin{equation*}
|\hat{\varphi}(\xi)| \leq \frac{C_{1}}{|\xi|^{n-1 / 2-\delta}} e^{-C_{2}|\xi|} \tag{2.16}
\end{equation*}
$$

for some $C_{1} \in \mathbf{R}, C_{2} \in \mathbf{R}$ from Stein [19, p. 132].
To prove (2.9), let $v(\xi)=h+|q(h \xi)|$ and $w(\xi)=|q(h \xi)|$. It suffices to prove that the operator $S$ with $S f=v^{1 / 2}\left(\hat{\varphi} * \frac{f}{w^{1 / 2}}\right)$ is bounded on $L^{2}$. Suppose that $S^{*}$ is the
adjoint of $S$, since $\left\|S^{*}\right\|_{L^{2} \rightarrow L^{2}}=\|S\|_{L^{2} \rightarrow L^{2}}$, it suffices to show that $S^{*}$ is bounded on $L^{2}$. We have

$$
\begin{align*}
\left\langle S^{*} f, g\right\rangle & =\langle f, S g\rangle \\
& =\int f(\xi) v(\xi)^{1 / 2} \overline{\left(\int \hat{\varphi}(\xi-\eta) \frac{g(\eta)}{w(\eta)^{1 / 2}} d \eta\right)} d \xi  \tag{2.17}\\
& =\int w(\eta)^{-1 / 2}\left(\int \overline{\hat{\varphi}(\xi-\eta)} v(\xi)^{1 / 2} f(\xi) d \xi\right) \overline{g(\eta)} d \eta
\end{align*}
$$

Thus

$$
S^{*} f(\xi)=w(\xi)^{-1 / 2} \int \overline{\hat{\varphi}(\eta-\xi)} v(\eta)^{1 / 2} f(\eta) d \eta
$$

Since

$$
\begin{aligned}
\left\|S^{*} f\right\|_{2}^{2} & = \\
& =\int\left(\int \overline{\hat{\varphi}(\eta-\xi)} v(\eta)^{1 / 2} f(\eta) d \eta\right) \cdot \overline{\left(\int \overline{\hat{\varphi}(\tilde{\eta}-\xi)} v(\tilde{\eta})^{1 / 2} f(\tilde{\eta}) d \tilde{\eta}\right)} w(\xi)^{-1} d \xi \\
& =\iiint \overline{\hat{\varphi}(\eta-\xi)} v(\eta)^{1 / 2} f(\eta) \hat{\varphi}(\tilde{\eta}-\xi) v(\tilde{\eta})^{1 / 2} \overline{f(\tilde{\eta})} w(\xi)^{-1} d \eta d \tilde{\eta} d \xi
\end{aligned}
$$

and

$$
\begin{aligned}
& \overline{\hat{\varphi}(\eta-\xi)} v(\eta)^{1 / 2} f(\eta) \hat{\varphi}(\tilde{\eta}-\xi) v(\tilde{\eta})^{1 / 2} \overline{f(\tilde{\eta})} \\
&=\left(\overline{\hat{\varphi}(\eta-\xi)} \hat{\varphi}(\tilde{\eta}-\xi)\left|\frac{\eta-\xi}{\tilde{\eta}-\xi}\right|^{1 / 2}\right)^{1 / 2} v(\eta)^{1 / 2} f(\eta) \\
& \cdot\left(\overline{\hat{\varphi}(\eta-\xi)} \hat{\varphi}(\tilde{\eta}-\xi)\left|\frac{\tilde{\eta}-\xi}{\eta-\xi}\right|^{1 / 2}\right)^{1 / 2} v(\tilde{\eta})^{1 / 2} \overline{f(\tilde{\eta})}
\end{aligned}
$$

Using $a \cdot b \leq \frac{1}{2}\left(a^{2}+b^{2}\right)$, we have

$$
\left\|S^{*} f\right\|_{2}^{2} \leq \iiint|\overline{\hat{\varphi}(\eta-\xi)}||\hat{\varphi}(\tilde{\eta}-\xi)|\left|\frac{\eta-\xi}{\tilde{\eta}-\xi}\right|^{1 / 2} v(\eta)|f(\eta)|^{2} w(\xi)^{-1} d \tilde{\eta} d \eta d \xi
$$

Now we integrate in the variable $\tilde{\eta}$, and use that

$$
\int|\hat{\varphi}(\tilde{\eta}-\xi)|\left|\frac{1}{\xi-\tilde{\eta}}\right|^{1 / 2} d \tilde{\eta} \leq C
$$

This implies that

$$
\begin{aligned}
\left\|S^{*} f\right\|_{2}^{2} & \leq C \iint|\overline{\hat{\varphi}(\eta-\xi)}||\eta-\xi|^{1 / 2} v(\eta)|f(\eta)|^{2} w(\xi)^{-1} d \eta d \xi \\
& \leq C \int\left(\int|\overline{\hat{\varphi}(\eta-\xi)}||\eta-\xi|^{1 / 2} v(\eta) w(\xi)^{-1} d \xi\right)|f(\eta)|^{2} d \eta
\end{aligned}
$$

It follows that

$$
\begin{equation*}
\left\|S^{*}\right\|_{L^{2} \rightarrow L^{2}} \leq C \max _{\eta}\left(\int \overline{\hat{\varphi}(\eta-\xi)}| | \eta-\left.\xi\right|^{1 / 2} v(\eta) w(\xi)^{-1} d \xi\right)^{1 / 2} \tag{2.18}
\end{equation*}
$$

We need to show the right side of $(2.18)$ is bounded by some finite number.
Since $|q(h \xi)| \approx|h \xi|^{2}$ when $|h \xi|>4$ and $|q(h \eta)| \leq|q(h \xi)|+|q(h(\eta-\xi))|+h^{2}|\xi| \mid \eta-$ $\xi \mid$, we have

$$
\begin{align*}
\left|\frac{v(\eta)}{w(\xi)}\right| & =\left|\frac{h+|q(h \eta)|}{|q(h \xi)|}\right| \\
& \left.\leq C \frac{h+|q(h \xi)|+|q(h(\eta-\xi))|+h^{2}|\xi||\eta-\xi|}{\mid q(h \xi)} \right\rvert\, \\
& \leq C\left(1+\frac{h+h^{2}|\eta-\xi|^{2}+h|\eta-\xi|+h^{2}|\xi||\eta-\xi|}{|q(h \xi)|}\right)  \tag{2.19}\\
& \leq C\left(1+\frac{h+h^{2}|\eta-\xi|^{2}+h|\eta-\xi|}{|q(h \xi)|}\right)
\end{align*}
$$

Now we consider the right side of (2.18),

$$
\begin{align*}
& \left|\int\right| \overline{\hat{\varphi}(\eta-\xi)}\left||\eta-\xi|^{1 / 2} v(\eta) w(\xi)^{-1} d \xi\right| \\
& \quad \leq C \int e^{-C_{2}|\eta-\xi|} \frac{1}{|\eta-\xi|^{n-1 / 2-\delta}}|\eta-\xi|^{1 / 2}\left(1+\frac{h+h^{2}|\eta-\xi|^{2}+h|\eta-\xi|}{|q(h \xi)|}\right) d \xi \\
& \quad \leq C \int e^{-C_{2}|\eta-\xi|} \frac{1}{|\eta-\xi|^{n-1-\delta}}\left(1+\frac{h+h^{2}|\eta-\xi|^{2}+h|\eta-\xi|}{|q(h \xi)|}\right) d \xi \\
& \quad \leq C \int e^{-C_{2}|\eta-\xi|} \frac{1}{|\eta-\xi|^{n-1-\delta}}\left(1+\frac{h}{|q(h \xi)|}\right) d \xi \\
& \quad \leq C+C \int e^{-C_{2}|\eta-\xi|} \frac{1}{|\eta-\xi|^{n-1-\delta}} \frac{h}{|q(h \xi)|} d \xi \tag{2.20}
\end{align*}
$$

For the integral

$$
\begin{align*}
& \int e^{-C_{2}|\eta-\xi|} \frac{1}{|\eta-\xi|^{n-1-\delta}} \frac{h}{|q(h \xi)|} d \xi \\
& \leq \int_{|\eta-\xi|>1} e^{-C_{2}|\eta-\xi|} \frac{1}{|\eta-\xi|^{n-1-\delta}} \frac{h}{|q(h \xi)|} d \xi \\
& \quad+\int_{B_{1}(\eta)} e^{-C_{2}|\eta-\xi|} \frac{1}{|\eta-\xi|^{n-1-\delta}} \frac{h}{|q(h \xi)|} d \xi  \tag{2.21}\\
& \leq \int_{|\eta-\xi|>1} e^{-C_{2}|\eta-\xi|} \frac{h}{|q(h \xi)|} d \xi+\int_{B_{1}(\eta)} \frac{1}{|\eta-\xi|^{n-1-\delta}} \frac{h}{|q(h \xi)|} d \xi \\
& =A+B
\end{align*}
$$

By estimate 2.2 in Lemma 2.1, $A$ is bounded by finite number. For $B$, by
estimate (2.1) in Lemma 2.1, we have

$$
\begin{align*}
B & \leq \sum_{k=0}^{\infty} \int_{B_{2-k}(\eta) \backslash B_{2}-k-1}(\eta) \\
& \leq C \sum_{k=0}^{\infty} 2^{k(n-1-\delta)} 2^{-k(n-1)}  \tag{2.22}\\
& \leq C \sum_{k=0}^{\infty} 2^{-k \delta} \\
& \leq C
\end{align*}
$$

Thus we have proven the estimate (2.9). The estimate (2.10) follows from (2.9) by duality. For the estimate 2.11 , since $1<\frac{h+|q(h \xi)|}{h}$, it follows that $\left\|\langle x\rangle^{-1 / 2-\delta} f\right\|_{L^{2}} \leq$ $h^{-1 / 2}\left\|\langle x\rangle^{-1 / 2-\delta} f\right\|_{X_{\mu, h}^{1 / 2}}$, which is bounded by $C h^{-1 / 2}\|f\|_{\dot{X}_{\mu, h}^{1 / 2}}$ from estimate 2.9 . The estimates (2.12) and (2.13) follow easily from the definition of the Bourgain type spaces.

## The semiclassical pseudodifferential operators

In our proof of existence of solutions, we need to use semiclassical pseudodifferential operators and we give the definition of these operators.

Definition 2.3. Let $0 \leq \sigma_{1}, \sigma_{2} \leq 1$ with $\sigma_{1}+\sigma_{2} \leq 1$. We define $S_{\sigma_{1}, \sigma_{2}}$ to be the space of all functions $a(x, \xi ; h)$ where $x, \xi \in \mathbf{R}^{n}$ and $0<h<h_{0}, h_{0}<1$, such that for each $h$, we have $a(x, \xi ; h) \in C^{\infty}\left(\mathbf{R}^{2 n}\right)$ and

$$
\partial_{x}^{\alpha} \partial_{\xi}^{\beta} a(x, \xi ; h) \leq C_{\alpha \beta} h^{-\sigma_{1}|\alpha|-\sigma_{2}|\beta|}
$$

for any multi-index $\alpha, \beta$.
Next we give more general symbol classes, which are defined by using an order function. This is similar to the symbol classes in R. Beals [1, p.3] but we give a semiclassical version of his definition.

Definition 2.4. Let $\varphi(x, \xi ; h): \mathbf{R}^{n} \times \mathbf{R}^{n} \times\left(0, h_{0}\right] \rightarrow[0, \infty)$, we say $\varphi$ is an order function if there exist $\sigma_{1} \geq 0, \sigma_{2} \geq 0, N_{1} \geq 0$ and $N_{2} \geq 0$ so that $\varphi$ satisfies

$$
|\varphi(x, \xi ; h)| \leq C\left\langle h^{-\sigma_{1}}(x-y)\right\rangle^{N_{1}}\left\langle h^{-\sigma_{2}}(\xi-\eta)\right\rangle^{N_{2}}|\varphi(y, \eta ; h)|
$$

for all $x, y, \xi, \eta \in \mathbf{R}^{n}$.
Definition 2.5. Let $\sigma_{1}, \sigma_{2} \in[0,1]$, and $\sigma_{1}+\sigma_{2} \leq 1$, we define the symbol class $S_{\sigma_{1}, \sigma_{2}}^{\varphi}\left(\mathbf{R}^{n}\right)$ as the space of all functions a $: \mathbf{R}^{n} \times \mathbf{R}^{n} \times\left(0, h_{0}\right] \rightarrow \mathbf{C}$ such that $a(\cdot ; h) \in$ $C^{\infty}\left(\mathbf{R}^{2 n}\right)$ and

$$
\left|\partial_{x}^{\alpha} \partial_{\xi}^{\beta} a(x, \xi ; h)\right| \leq C_{\alpha, \beta} h^{-|\alpha| \sigma_{1}-|\beta| \sigma_{2}} \varphi(x, \xi ; h) .
$$

When the order function $\varphi=1$, we use $S_{\sigma_{1}, \sigma_{2}}$ for symbols in place of $S_{\sigma_{1}, \sigma_{2}}^{1}$ in order to avoid confusion with other common notations. Given a symbol $a$ in one of the classes defined above, there is a natural way to define an operator $A=\mathrm{Op}_{h}(a)$ : $\mathcal{S}\left(\mathbf{R}^{n}\right) \rightarrow \mathcal{S}^{\prime}\left(\mathbf{R}^{n}\right)$. This is the subject of the next definition.

Definition 2.6. For $a \in S_{\sigma_{1}, \sigma_{2}}^{\varphi}$, we define the operator $A=O p_{h}(a)=a(x, h D)$ by

$$
A f(x)=(2 \pi)^{-n} \int_{\mathbf{R}^{n}} e^{i x \cdot \xi} a(x, h \xi ; h) \hat{f}(\xi) d \xi
$$

where $\hat{f}$ is the Fourier transform of $f$. The function $a(x, \xi ; h)$ is called the symbol of the operator $A$.

Lemma 2.7. Let $b \in \mathbf{R}$, then the function $\phi(\xi)=\left(h^{2(1-\sigma)}+|q(\xi)|^{2}\right)^{b / 2}$ is an order function and a symbol in $S_{0,1-\sigma}^{\phi}$.

Proof. We know $|q(\xi)|=|q(\xi-\eta+\eta)| \leq|q(\xi-\eta)|+|q(\eta)|+2|\xi-\eta||\eta|$. Thus

$$
\begin{align*}
& \frac{\left(h^{2(1-\sigma)}+|q(\xi)|^{2}\right)^{b / 2}}{\left(h^{2(1-\sigma)}+|q(\eta)|^{2}\right)^{b / 2}} \leq \leq\left(1+\frac{(|q(\xi-\eta)|+2|\xi-\eta||\eta|)^{2}}{h^{2(1-\sigma)}+|q(\eta)|^{2}}\right)^{b / 2} \\
& \leq\left(1+\frac{\left(|\xi-\eta|^{2}+2|\mu \cdot(\xi-\eta)|+2|\xi-\eta||\eta|\right)^{2}}{h^{2(1-\sigma)}+|q(\eta)|^{2}}\right)^{b / 2} \\
& \leq\left(1+\left(\left|h^{-2(1-\sigma)}\right| \xi-\left.\eta\right|^{4}+2 h^{-2(1-\sigma)}|\xi-\eta|^{2}\right.\right.  \tag{2.23}\\
&\left.\left.\quad+\frac{(2|\xi-\eta||\eta|)^{2}}{h^{2(1-\sigma)}+|q(\eta)|^{2}}\right)\right)^{b / 2} \\
& \leq C\left\langle h^{-(1-\sigma)}(\xi-\eta)\right\rangle^{2 b} .
\end{align*}
$$

which gives us $\phi$ is an order function.
To prove $\phi$ is a symbol, we need to estimate the norm of derivatives. We estimate the first order partial derivatives as follows

$$
\begin{align*}
\left|\frac{\partial}{\partial \xi_{j}} \phi(\xi)\right| & \left.=\left\lvert\, \frac{b}{2}\left(h^{2(1-\sigma)}+|q(\xi)|^{2}\right)^{b / 2-1} \frac{\partial}{\partial \xi_{j}}(q(\xi) q \bar{\xi})\right.\right) \mid \\
& =\left|\frac{b}{2}\right|\left|\left(h^{2(1-\sigma)}+|q(\xi)|^{2}\right)^{-1}\left(\left(2 \xi_{j}+2 \mu_{j}\right) \overline{q(\xi)}+\overline{\left(2 \xi_{j}+2 \mu_{j}\right)} q(\xi)\right)\right||\phi(\xi)| \\
& \leq C\left|\left(h^{2(1-\sigma)}+|q(\xi)|^{2}\right)^{-1 / 2}\right| 2 \xi_{j}+2 \mu_{j}| ||\phi(\xi)| \\
& \leq C h^{-(1-\sigma)}|\phi(\xi)| \tag{2.24}
\end{align*}
$$

Similarly, we can calculate higher partial derivatives, and get $\phi \in S_{0,1-\sigma}^{\phi}$.
Lemma 2.8. If $a \in S_{\sigma_{1}, \sigma_{2}}^{\varphi}, b \in S_{\sigma_{1}, \sigma_{2}}^{\psi}$, then $a b \in S_{\sigma_{1}, \sigma_{2}}^{\varphi \psi}$.

Proof. By the product rule,

$$
\begin{align*}
\left|\partial_{x}^{\alpha} \partial_{\xi}^{\beta} a(x, \xi ; h) b(x, \xi ; h)\right| & =\left|\sum_{\substack{\alpha_{1}+\alpha_{2}=\alpha \\
\beta_{1}+\beta_{2}=\beta}} C\left(\alpha_{1}, \beta_{1}, \alpha_{2}, \beta_{2}\right) \partial_{x}^{\alpha_{1}} \partial_{\xi}^{\beta_{1}} a \partial_{x}^{\alpha_{2}} \partial_{\xi}^{\beta_{2}} b\right| \\
& \leq\left|\sum_{\substack{\alpha_{1}+\alpha_{2}=\alpha \\
\beta_{1}+\beta_{2}=\beta}} C\left(\alpha_{1}, \beta_{1}, \alpha_{2}, \beta_{2}\right) h^{-\left|\alpha_{1}\right| \sigma_{1}-\left|\beta_{1}\right| \sigma_{2}} \varphi h^{-\left|\alpha_{2}\right| \sigma_{1}-\left|\beta_{2}\right| \sigma_{2}} \phi\right| \\
& \leq C(\alpha, \beta) h^{-|\alpha| \sigma_{1}-|\beta| \sigma_{2}} \varphi \phi \tag{2.25}
\end{align*}
$$

Thus $a b \in S_{\sigma_{1}, \sigma_{2}}^{\varphi \psi}$.
The lemma below shows that if $A$ and $B$ are semiclassical pseudodifferential operators, then the composition $A \circ B$ is a semiclassical pseudodifferential operator and gives information about the symbol.

Lemma 2.9. Let $a \in S_{\sigma_{1}, \sigma_{2}}^{\varphi}$ and $b \in S_{\tilde{\sigma}_{1}, \tilde{\sigma}_{2}}^{\psi}$, with $\sigma_{2}+\tilde{\sigma}_{1} \leq 1, c(x, \xi)$ be the symbol of operator $A \circ B$ with $A=O p_{h}(a)$ and $B=O p_{h}(b)$, then $c \in S_{\max \left(\sigma_{1}, \tilde{\sigma}_{1}\right), \max \left(\sigma_{2}, \tilde{\sigma}_{2}\right)}^{\varphi \psi}$ and for each $M=1,2, \ldots$ we have

$$
\begin{equation*}
c(x, \xi)=\sum_{|\alpha|<M} \frac{h^{|\alpha|}}{\alpha!} \partial_{\xi}^{\alpha} a(x, \xi) D_{x}^{\alpha} b(x, \xi)+h^{M\left(1-\sigma_{2}-\tilde{\sigma}_{1}\right)} S_{\sigma_{1}, \tilde{\sigma}_{2}}^{\varphi \psi} \tag{2.26}
\end{equation*}
$$

Proof. Our argument follows the proof in Stein [20, page 320-323], with obvious changes to handle the parameter $h$ and the order functions.

We first assume that $a$ and $b$ have compact support in the $x$ and $\xi$ variables. By definition

$$
\begin{aligned}
& A f(x)=(2 \pi)^{-n} \int a(x, h \eta) e^{i \eta(x-y)} f(y) d y d \eta \\
& B f(y)=(2 \pi)^{-n} \int b(y, h \xi) e^{i \xi(y-z)} f(z) d z d \xi
\end{aligned}
$$

So we have

$$
A B f(x)=(2 \pi)^{-2 n} \int a(x, h \eta) e^{i \eta \cdot(x-y)} b(y, h \xi) e^{i \xi \cdot(y-z)} f(z) d z d \xi d y d \eta
$$

Since $e^{i \eta(x-y)} e^{i \xi(y-z)}=e^{i \xi(x-z)} e^{i(\eta-\xi)(x-y)}$, we can write

$$
A B f(x)=(2 \pi)^{-n} \int e^{i \xi \cdot(x-z)} c(x, h \xi) f(z) d z d \xi
$$

where

$$
\begin{aligned}
c(x, h \xi) & =(2 \pi)^{-n} \int e^{i(\eta-\xi) \cdot(x-y)} a(x, h \eta) b(y, h \xi) d y d \eta \\
& =(2 \pi)^{-n} \int e^{i \eta \cdot(x-y)} a(x, h(\eta+\xi)) b(y, h \xi) d y d \eta .
\end{aligned}
$$

If we replace $h \xi$ with $\xi$ in the definition of $c(x, h \xi)$, we have

$$
\begin{equation*}
c(x, \xi)=(2 \pi)^{-n} \int e^{i \eta \cdot(x-y)} a(x, \xi+h \eta) b(y, \xi) d y d \eta \tag{2.27}
\end{equation*}
$$

or equivalently

$$
\begin{equation*}
c(x, \xi)=(2 \pi)^{-n} \int e^{i \eta \cdot x} a(x, \xi+h \eta) \hat{b}(\eta, \xi) d \eta \tag{2.28}
\end{equation*}
$$

where $\hat{b}$ denotes the Fourier transform of $b$ in the first variable.
We choose an arbitrary point $x_{0}$ in $\mathbf{R}^{n}$ and we will compute $c(x, \xi)$ for $x$ with $\left|x-x_{0}\right|<\frac{1}{2} h^{\tilde{\sigma}_{1}}$. Let $\chi(x) \in C_{c}^{\infty}\left(\mathbf{R}^{n}\right)$, where $\chi(x)=1$, when $x \in B\left(x_{0}, h^{\tilde{\sigma}_{1}}\right)$, and $\chi(x)=0$, when $x$ outside of $B\left(x_{0}, 2 h^{\tilde{\sigma}_{1}}\right)$. Define $b_{0}(x, \xi)=b(x, \xi) \chi(x)$, and $b_{\infty}(x, \xi)=b(x, \xi)(1-\chi(x))$, define $c_{0}(x, \xi)$ and $c_{\infty}(x, \xi)$ as

$$
\begin{aligned}
c_{0}(x, \xi) & =(2 \pi)^{-n} \int e^{i \eta \cdot(x-y)} a(x, \xi+h \eta) b_{0}(y, \xi) d y d \eta \\
c_{\infty}(x, \xi) & =(2 \pi)^{-n} \int e^{i \eta \cdot(x-y)} a(x, \xi+h \eta) b_{\infty}(y, \xi) d y d \eta
\end{aligned}
$$

Thus $c(x, \xi)=c_{0}(x, \xi)+c_{\infty}(x, \xi)$.
We first consider $c_{0}(x, \xi)$. Replacing $b$ by $b_{0}$ in (2.28), we have

$$
c_{0}(x, \xi)=(2 \pi)^{-n} \int e^{i \eta \cdot x} a(x, \xi+h \eta) \hat{b}_{0}(\eta, \xi) d \eta
$$

Apply Taylor's formula to $a(x, \xi+h \eta)$, to obtain

$$
a(x, \xi+h \eta)=\sum_{|\alpha|<M} \frac{1}{\alpha!} \partial_{\xi}^{\alpha} a(x, \xi)(h \eta)^{\alpha}+R_{M}(x, \xi, \eta, h)
$$

For each $\alpha$, we obtain

$$
(2 \pi)^{-n} \int e^{i \eta \cdot x} \frac{1}{\alpha!} \partial_{\xi}^{\alpha} a(x, \xi)(h \eta)^{\alpha} \hat{b}_{0}(\eta, \xi) d y d \eta=\frac{h^{|\alpha|}}{\alpha!} \partial_{\xi}^{\alpha} a(x, \xi) D_{x}^{\alpha} b_{0}(x, \xi) .
$$

Then for the remainder part, we know

$$
\begin{aligned}
\left|R_{M}(x, \xi, \eta, h)\right| & \leq C_{M} \sup _{t \in[0,1],|\alpha|=M}\left\{\left|\partial_{\xi^{\prime}}^{\alpha} a\left(x, \xi^{\prime}\right)\right|: \xi^{\prime}=\xi+t h \eta\right\}|h \eta|^{M} \\
& \leq C_{M} h^{M\left(1-\sigma_{2}\right)}|\eta|^{M} \sup _{t \in[0,1],|\alpha|=M}\left\{\left|\varphi\left(x, \xi^{\prime}\right)\right|: \xi^{\prime}=\xi+t h \eta\right\} \\
& \leq C_{M} h^{M\left(1-\sigma_{2}\right)}|\eta|^{M}\left\langle h^{1-\sigma_{2}} \eta\right\rangle^{N_{2}} \varphi(x, \xi)
\end{aligned}
$$

and since $b_{0}$ is in the symbol class $S_{\tilde{\sigma}_{1}, \tilde{\sigma}_{2}}^{\phi}$, we have

$$
\left|\hat{b}_{0}(\eta, \xi)\right|=\left|\int \frac{e^{-i y \eta}}{\left\langle h^{\tilde{\sigma}_{1}} \eta\right\rangle^{2 M_{1}}}\left(1-h^{2 \tilde{\sigma}_{1}} \Delta_{y}\right)^{M_{1}} b_{0}(y, \xi) d y\right| \leq C_{M_{1}} \frac{h^{\tilde{\sigma}_{1} n}}{\left\langle h^{\tilde{\sigma}_{1}} \eta\right\rangle^{2 M_{1}}}|\psi(x, \xi)| .
$$

So the remainder integral will be bounded by

$$
C_{M, M_{1}} h^{M\left(1-\sigma_{2}\right)}|\varphi(x, \xi) \| \psi(x, \xi)| \int \frac{|\eta|^{M}\left\langle h^{1-\sigma_{2}} \eta\right\rangle^{N_{2}} h^{\tilde{\sigma}_{1} n}}{\left\langle h^{\tilde{\sigma}_{1}} \eta\right\rangle^{2 M_{1}}} d \eta .
$$

If we choose $M_{1}$ large enough so the integral in this estimate is integrable, this bound will be less then

$$
C_{M} h^{M\left(1-\sigma_{2}-\tilde{\sigma}_{1}\right)}|\varphi(x, \xi)||\psi(x, \xi)|
$$

Further, if we apply $\partial_{x}^{\alpha} \partial_{\xi}^{\beta}$ to the remainder, the remainder will be bounded by $C_{M} h^{M\left(1-\sigma_{2}-\tilde{\sigma}_{1}\right)} h^{-\sigma_{1}|\alpha|-\tilde{\sigma}_{2}|\beta|}|\varphi(x, \xi)||\psi(x, \xi)|$, and we can conclude the remainder is in $h^{M\left(1-\sigma_{2}-\tilde{\sigma}_{1}\right)} S_{\sigma_{1}, \tilde{\sigma}_{2}}^{\varphi \psi}$.

Above all we have

$$
c_{0}(x, \xi)=\sum_{|\alpha|<M} \frac{h^{|\alpha|}}{\alpha!} \partial_{\xi}^{\alpha} a(x, \xi) D_{x}^{\alpha} b(x, \xi)+h^{M\left(1-\sigma_{2}-\tilde{\sigma}_{1}\right)} S_{\sigma_{1}, \tilde{\sigma}_{2}}^{\varphi \psi}
$$

Next we deal with $c_{\infty}(x, \xi)$. Since

$$
\begin{gathered}
(-\Delta)_{\eta}^{M_{1}} e^{i \eta \cdot(x-y)}=(-1)^{M_{1}}|x-y|^{2 M_{1}} e^{i \eta(x-y)} \\
\left(1-h^{2 \tilde{\sigma}_{1}} \Delta_{y}\right)^{M_{2}} e^{i \eta \cdot(x-y)}=\left\langle h^{\tilde{\sigma}_{1}} \eta\right\rangle^{2 M_{2}} e^{i \eta(x-y)} \\
|x-y|>\frac{1}{2} h^{\tilde{\sigma}_{1}}
\end{gathered}
$$

then,

$$
\begin{aligned}
\left|c_{\infty}(x, \xi)\right|= & \left|\int e^{i \eta \cdot(x-y)} \frac{\Delta_{\eta}^{M_{1}} a(x, \xi+h \eta)}{\left\langle h^{\tilde{\sigma}_{1}} \eta\right\rangle^{2 M_{2}}}\left(1-h^{2 \tilde{\sigma}_{1}} \Delta_{y}\right)^{M_{2}}\left\{\frac{b_{\infty}(y, \xi)}{|x-y|^{2 M_{1}}}\right\} d y d \eta\right| \\
\leq & C_{M_{1}, M_{2}} \left\lvert\, \int \frac{h^{2 M_{1}\left(1-\sigma_{2}-\tilde{\sigma}_{1}\right.}|\varphi(x, \xi+h \eta)|}{\left\langle h^{\tilde{\sigma}_{1}} \eta\right\rangle^{M_{2}}}\left(1-h^{2 \tilde{\sigma}_{1}} \Delta_{y}\right)^{M_{2}}\right. \\
& \left.\times\left\{\frac{b_{\infty}(y, \xi)}{\left(h^{-\tilde{\sigma}_{1}}|x-y|\right)^{2 M_{1}}}\right\} d y d \eta \right\rvert\, \\
\leq & C_{M_{1}, M_{2}} \left\lvert\, \int \frac{h^{2 M_{1}\left(1-\sigma_{2}-\tilde{\sigma}_{1}\right)}\left\langle h^{-\sigma_{2}} h \eta\right\rangle^{N_{2}}|\varphi(x, \xi)|}{\left\langle h^{\tilde{\sigma}_{1}} \eta\right\rangle^{2 M_{2}}}\right. \\
& \left.\times \frac{|\psi(y, \xi)|}{\left(h^{-\tilde{\sigma}_{1}}|x-y|\right)^{2 M_{1}}} d y d \eta \right\rvert\, \\
\leq & \left.C_{M_{1}, M_{2}} h^{2 M_{1}\left(1-\sigma_{2}-\tilde{\sigma}_{1}\right)} \int \frac{\left\langle h^{-\sigma_{2}} h \eta\right\rangle^{N_{2}}|\varphi(x, \xi)|}{\left\langle h \tilde{\sigma}_{1}\right.} \eta\right\rangle^{2_{2}} \\
& \times \frac{\left\langle h^{-\tilde{\sigma}_{1}}(x-y)\right\rangle^{N_{1}}|\psi(x, \xi)|}{\left(h^{-\tilde{\sigma}_{1}}|x-y|\right)^{2 M_{1}}} d y d \eta
\end{aligned}
$$

If we choose $M_{1}, M_{2}$ large enough, we get this integral is less then

$$
\begin{aligned}
& C_{M_{1}} h^{2 M_{1}\left(1-\sigma_{2}-\tilde{\sigma}_{1}\right)}|\varphi(x, \xi)||\psi(x, \xi)| \\
& \leq C_{M} h^{M\left(1-\sigma_{2}-\tilde{\sigma}_{1}\right)}|\varphi(x, \xi)||\psi(x, \xi)|
\end{aligned}
$$

Thus, we know that

$$
c(x, \xi)=\sum_{|\alpha|<M} \frac{h^{|\alpha|}}{\alpha!} \partial_{\xi}^{\alpha} a(x, \xi) D_{x}^{\alpha} b(x, \xi)+h^{M\left(1-\sigma_{2}-\tilde{\sigma}_{1}\right)} S_{\sigma_{1}, \tilde{\sigma}_{2}}^{\varphi \psi}
$$

Finally, we consider the case when $a$ and $b$ do not have compact support. We choose a function $\gamma(x, \xi) \in C_{c}^{\infty}\left(\mathbf{R}^{n} \times \mathbf{R}^{n}\right)$, with $\gamma(0)=1$, then we can use same method to get 2.27) and 2.28 by replacing $a$ and $b$ with $a_{\tau}$ and $b_{\tau}$ respectively, where $a_{\tau}(x, \xi)=a(x, \xi) \cdot \gamma(\tau x, \tau \xi)$ and $b_{\tau}(x, \xi)=b(x, \xi) \cdot \gamma(\tau x, \tau \xi)$. Note that $a_{\tau}$ and $b_{\tau}$ are in the same symbol class as $a$ and $b$ respectively, $0<\tau<1$. We let $\tau \rightarrow 0$ to get the symbol $c$. Let $a_{\tau}(x, \xi)=a(x, \xi) \cdot \gamma(\tau x, \tau \xi)$ and $b_{\tau}(x, \xi)=b(x, \xi) \cdot \gamma(\tau x, \tau \xi)$, and define

$$
C_{\tau}=A_{\tau} \circ B_{\tau}
$$

We have proven that $c_{\tau}$, which is the symbol of $C_{\tau}$, satisfies the formula (2.26), with $c_{\tau}, a_{\tau}$ and $b_{\tau}$ replacing $c, a$ and $b$ respectively, uniformly in $\tau$. And what we did above show that $c_{\tau}$ converges pointwise to some limit $c$. By the continuity properties (see Stein [20, pp. 232-233]), we have $C=A \circ B$. Thus, this theorem is proved.

We use $L_{r}^{2}$ to denote the weighted $L^{2}$ space with the norm $\|f\|_{L_{r}^{2}}=\left\|\langle\cdot\rangle^{r} f(\cdot)\right\|_{2}$. We need use this space in the following Proposition.

Proposition 2.10. Let $a \in S_{\sigma_{1}, \sigma_{2}}$ with $\sigma_{1}, \sigma_{2} \in[0,1], \sigma_{1}+\sigma_{2} \leq 1$, $A=O p_{h}(a)$ then
(1) $A$ is bounded on $L^{2}$.
(2) $A$ is bounded on $X_{\mu, h, \sigma}^{b}$, for any $1>\sigma \geq \sigma_{1}, b \in(-1,1)$.
(3) Suppose $A$ is bounded and invertible bounded in $L^{2}$, then

$$
\left\|A\langle x\rangle^{r} f\right\|_{L^{2}} \approx\left\|\langle x\rangle^{r} A f\right\|_{L^{2}}
$$

for any real number $r$.
Proof. For part (1), If $\sigma_{1}=\sigma_{2}$, the proof can be found in [7]. If $\sigma_{1} \neq \sigma_{2}$, since

$$
\|a(x, h D)\|_{L^{2} \rightarrow L^{2}}=\left\|a\left(s x, s^{-1} h D\right)\right\|_{L^{2} \rightarrow L^{2}}
$$

for $s \in \mathbf{R}$. So we can choose $s=h^{\frac{\sigma_{1}-\sigma_{2}}{2}}$, so that $a\left(s x, s^{-1} h D\right) \in S_{\frac{\sigma_{1}+\sigma_{2}}{2}, \frac{\sigma_{1}+\sigma_{2}}{2}}^{1}$, which is an operator bounded on $L^{2}$. Thus $a(x, h D)$ is bounded on $L^{2}$.

For (2), we let $\varphi(\xi)=\left(h^{2(1-\sigma)}+|q(\xi)|^{2}\right)^{b / 2} \in S_{0,1-\sigma}^{\varphi}$ which is from Lemma 2.7. The estimate (2) is equivalent to proving that $\varphi(h D) \circ a(x, h D) \circ \varphi(h D)^{-1}$ is bounded on $L^{2}$. Now we consider two operators $\varphi(h D)$ and $a(x, h D) \varphi(h D)^{-1}$. Since $(1-\sigma)+\sigma_{1} \leq 1$, then by using Lemma 2.9, we know that the composition of these two operators is a semiclassical pseudodifferential operator with symbol in $S_{\sigma, 1-\sigma}$, and this operator is bounded on $L^{2}$ by part (1) of this Proposition.

Finally for part (3), by Salo [18, Propostion 2.2], we know that the operator norm $\|A\|_{L_{r}^{2} \rightarrow L_{r}^{2}} \leq C$ for any $r$. Thus

$$
\begin{align*}
\left\|A\langle x\rangle^{r} f\right\|_{L^{2}} & =\left\|A\langle x\rangle^{r} A^{-1}\langle x\rangle^{-r}\langle x\rangle^{r} A f\right\|_{L^{2}} \\
& \leq C\left\|A^{-1}\langle x\rangle^{-r}\langle x\rangle^{r} A f\right\|_{L^{2}}  \tag{2.29}\\
& \leq C\left\|\langle x\rangle^{-r}\langle x\rangle^{r} A f\right\|_{L^{2}} \\
& \leq C\left\|\langle x\rangle^{r} A f\right\|_{L^{2}} .
\end{align*}
$$

Similarly, we obtain $\left\|\langle x\rangle^{r} A f\right\|_{L^{2}} \leq C\left\|A\langle x\rangle^{r} f\right\|_{L^{2}}$.

## Chapter 3 Existence of CGO solution in $\dot{X}_{\mu, h}^{1 / 2}$ spaces

We need some notation before we come to the main topic. Let $\Delta=\sum D_{j} \cdot D_{j}$, $\Delta_{\zeta}=\Delta+2 \zeta \cdot D$ and $D_{\zeta}=D+\zeta$, where $D=\left(D_{1}, D_{2}, \cdots, D_{n}\right)$.

For each $\zeta \in \mathbf{C}^{n}$ with $\zeta \cdot \zeta=0$, we want to find CGO solutions of the equation $H_{W, p} u=0$ with $u=e^{i x \cdot \zeta}(1+v)$. Substituting $u$ in $H_{W, p} u=0$, we obtain that $v$ solves the equation

$$
\left(\Delta_{\zeta}+2 W \cdot D_{\zeta}+(W \cdot W+D \cdot W+p)\right) v=-(2 W \cdot \zeta+W \cdot W+D \cdot W+p)
$$

Let $G=W \cdot W+D \cdot W+p$ and $f=-(2 W \cdot \zeta+G)$, we have

$$
\begin{equation*}
\left(\Delta_{\zeta}+2 W \cdot D_{\zeta}+G\right) v=f \tag{3.1}
\end{equation*}
$$

Let $\mu=h \zeta$ with $|\mu|=\sqrt{2}, h=\sqrt{2} /|\zeta|$. After multiplying $h^{2}$ on both sides of (3.1), we obtain an equivalent equation

$$
\begin{equation*}
\left(Q(h D)+2 h W \cdot(h D+\mu)+h^{2} G\right) v=h^{2} f \tag{3.2}
\end{equation*}
$$

where $Q(h D)$ is the semiclassical pseudodifferential operator with symbol $q(\xi)=$ $\xi \cdot \xi+2 \mu \cdot \xi$.

Now we assume $W \in C^{\lambda}$ for some $\lambda>1 / 2, m(x)$ is a standard mollifier, $m_{\kappa}(x)=$ $\frac{1}{\kappa^{n}} m\left(\frac{x}{\kappa}\right)$, and

$$
\begin{align*}
& W^{\sharp}=W * m_{\kappa} \\
& W^{b}=W-W^{\sharp} \\
& \left\|W^{b}\right\|_{\infty} \leq C h^{\sigma_{0} \lambda}  \tag{3.3}\\
& \left\|\partial_{x}^{\alpha} W^{\sharp}\right\|_{\infty} \leq C h^{\sigma_{0}(\lambda-|\alpha|)}, \quad|\alpha| \geq 1
\end{align*}
$$

where $\kappa=h^{-\sigma_{0}}$ and $\sigma_{0} \in(0,1 / 2)$. Then our equation (3.2) becomes

$$
\begin{equation*}
\left(Q+2 h W^{\sharp} \cdot(h D+\mu)+2 h W^{b} \cdot(h D+\mu)+h^{2} G\right) v=h^{2} f . \tag{3.4}
\end{equation*}
$$

The reason we decompose $W$ is that the term $Q+2 h W^{\sharp} \cdot(h D+\mu)$ in (3.4) is the main term compared with the remaining terms. This will be clear after we finish the following Lemma 3.3 and Theorem 3.4.

First we present a lemma from Haberman and Tataru [10, Lemma 2.1],
Lemma 3.1. Let $w_{1}$, $w_{2}$ be nonnegative weights, and $\phi$ be a rapidly decreasing function, then

$$
\|\phi * v\|_{L_{w_{2}}^{2}} \leq C \min \left\{\sup _{\xi} \sqrt{J(\xi, \eta) d \eta}, \sup _{\eta} \sqrt{J(\eta, \xi) d \xi}\right\}\|v\|_{L_{w_{1}}^{2}}
$$

where

$$
J(\xi, \eta)=|\phi(\xi-\eta)| \frac{w_{2}(\xi)}{w_{1}(\eta)}
$$

For symbols with some special properties, the following Lemma gives that the pseudodifferential operator is bounded on the Bourgain space $\dot{X}_{\mu, h}^{1 / 2}$.

Lemma 3.2. Suppose that $\sigma_{0}+\sigma \leq 1$. If we assume symbol $a \in S_{\sigma_{0}, \sigma}$ and $(a-1)\langle x\rangle \in$ $S_{\sigma_{0}, \sigma}$, then for $h$ small enough, we have

$$
\begin{equation*}
\|A\|_{\dot{X}_{\mu, h}^{1 / 2}, \dot{X}_{\mu, h}^{1 / 2}} \leq C h^{-\sigma_{0} / 2} \tag{3.5}
\end{equation*}
$$

If in addition $|a|$ is bounded away from 0 and $\sigma_{0}+\sigma<1$, then we have the same estimate for the operator $A^{-1}$,

$$
\begin{equation*}
\left\|A^{-1}\right\|_{\dot{X}_{\mu, h}^{1 / 2}, \dot{X}_{\mu, h}^{1 / 2}} \leq C h^{-\sigma_{0} / 2} \tag{3.6}
\end{equation*}
$$

Proof. For any $u$, since $A=(A-I)+I$ and estimate (2.12) and (2.14) in Proposition (2.2), we have

$$
\begin{align*}
\|A u\|_{\dot{X}_{\mu, h}^{1 / 2}} & \leq\|(A-I) u\|_{\dot{X}_{\mu, h}^{1 / 2}}+\|I u\|_{\dot{X}_{\mu, h}^{1 / 2}} \\
& \leq\left\|(A-I)\langle x\rangle\langle x\rangle^{-1} u\right\|_{X_{\mu, h, \sigma_{0}}^{1 / 2}}+\|u\|_{\dot{X}_{\mu, h}^{1 / 2}} \\
& \leq C\left\|\langle x\rangle^{-1} u\right\|_{X_{\mu, h, \sigma_{0}}^{1 / 2}}+\|u\|_{\dot{X}_{\mu, h}^{1 / 2}}  \tag{3.7}\\
& \leq C h^{-\sigma_{0} / 2}\left\|\langle x\rangle^{-1} u\right\|_{X_{\mu, h}^{1 / 2}}+\|u\|_{\dot{X}_{\mu, h}^{1 / 2}}
\end{align*}
$$

By estimate 2.9) in Proposition 2.2, we obtain

$$
\begin{align*}
& \leq C h^{-\sigma_{0} / 2}\|u\|_{\dot{X}_{\mu, h}^{1 / 2}}+\|u\|_{\dot{X}_{\mu, h}^{1 / 2}} \\
& \leq C h^{-\sigma_{0} / 2}\|u\|_{\dot{X}_{\mu, h}^{1 / 2}} \tag{3.8}
\end{align*}
$$

Thus we have proven estimate (3.5). For estimate (3.6), we consider an approximate right-inverse of $A$ of the form

$$
I+\sum_{j=0}^{N} h^{j} B_{j}
$$

with $B_{j}=O p_{h}\left(b_{j}\right)$ and $b_{j} \in S_{\sigma_{0}, \sigma}^{\langle x\rangle^{-1}}$. In fact, $(a-1) \in S_{\sigma_{0}, \sigma}^{\langle x\rangle)^{-1}}$. And

$$
\begin{equation*}
A \circ\left(I+\sum_{j=0}^{N} h^{j} B_{j}\right)=I+h^{N\left(1-\sigma_{0}-\sigma\right)} O p_{h}\left(b_{N+1}\right), \tag{3.9}
\end{equation*}
$$

with $b_{N+1} \in S_{\sigma_{0}, \sigma}^{<x>^{-1}}$. Then for $N$ large enough, we have that $A$ has a right inverse on $\dot{X}_{\mu, h}^{1 / 2}$, and by a similar argument, we can find a left inverse, so $A$ is invertible.

To find $b_{j}$, we use the symbol calculus to formally compute the symbol of the composite operator

$$
\begin{align*}
A \circ\left(I+\sum_{j=0}^{N} h^{j} B_{j}\right)= & O p_{h}\left(a\left(1+b_{0}\right)\right. \\
& +h \sum_{|\alpha|=1} \partial_{\xi}^{\alpha} a D_{x}^{\alpha} b_{0}+h a b_{1} \\
& +h^{2} \sum_{|\alpha|=2} \frac{1}{\alpha!} \partial_{\xi}^{\alpha} a D_{x}^{\alpha} b_{0}+\sum_{|\alpha|=1} \partial_{\xi}^{\alpha} a D_{x}^{\alpha} b_{0}+h^{2} a b_{2}  \tag{3.10}\\
& +\cdots \\
& +h^{N} \sum_{|\alpha|=N} \frac{1}{\alpha!} \partial_{\xi}^{\alpha} a D_{x}^{\alpha} b_{0}+\cdots+h^{N} a b_{N} \\
& \left.+h^{N\left(1-\sigma_{0}-\sigma\right)} O p_{h}\left(b_{N+1}\right)\right),
\end{align*}
$$

where $b_{N+1} \in S_{\sigma_{0}, \sigma}^{\langle x\rangle^{-1}}$ by Lemma 2.9. Solving for $b_{j}$ to make all but the remainder term vanish gives

$$
\begin{align*}
& b_{0}=\frac{1}{a}-1 \\
& b_{1}=-\frac{1}{a}\left(\sum_{|\alpha|=1} \partial_{\xi}^{\alpha} a D_{x}^{\alpha} b_{0}\right) \\
& b_{2}=-\frac{1}{a}\left(\sum_{|\alpha|=2} \frac{1}{\alpha!} \partial_{\xi}^{\alpha} a D_{x}^{\alpha} b_{0}+\sum_{|\alpha|=1} \partial_{\xi}^{\alpha} a D_{x}^{\alpha} b_{1}\right)  \tag{3.11}\\
& \vdots \\
& b_{N}=-\frac{1}{a}\left(\sum_{|\alpha|=N} \frac{1}{\alpha!} \partial_{\xi}^{\alpha} a D_{x}^{\alpha} b_{0}+\cdots+\sum_{|\alpha|=1} \partial_{\xi}^{\alpha} a D_{x}^{\alpha} b_{N-1}\right)
\end{align*}
$$

Since $I+\sum_{j=1}^{N} h^{j} B_{j}$ and $\left(I+h^{N\left(1-\sigma_{0}-\sigma\right)} O p_{h}\left(b_{N+1}\right)\right)^{-1}$ is bounded on $\dot{X}_{\mu, h}^{1 / 2}$ with norm $C h^{-\sigma_{0} / 2}$ and $C$ respectively, thus $A^{-1}=\left(I+\sum_{j=1}^{N} h^{j} B_{j}\right) \circ\left(I+h^{N\left(1-\sigma_{0}-\sigma\right)} O p_{h}\left(b_{N+1}\right)\right)^{-1}$ is bounded and $\left\|A^{-1}\right\|_{\dot{X}_{\mu, h}^{1 / 2}, \dot{X}_{\mu, h}^{1 / 2}} \leq C h^{-\sigma_{0} / 2}$.

Then we prove the following lemma.
Lemma 3.3. Suppose that $h$ is small, $\delta>0$ is small and $W \in C_{c}^{\lambda}\left(\mathbf{R}^{n}\right)$ with $\lambda>1 / 2$. If $g$ satisfies $\langle x\rangle^{1 / 2+\delta} g \in \dot{X}_{\mu, h}^{-1 / 2}$, then we have a solution $v$ to the equation

$$
\begin{equation*}
\left(Q(h D)+2 h W^{\sharp} \cdot(h D+\mu)\right) v=g \tag{3.12}
\end{equation*}
$$

with

$$
\begin{equation*}
\|v\|_{\dot{X}_{\mu, h}^{1 / 2}} \leq C h^{-\sigma_{0}}\left\|\langle x\rangle^{1 / 2+\delta} g\right\|_{X_{\mu, h}^{-1 / 2}} . \tag{3.13}
\end{equation*}
$$

Furthermore, there is only one solution to (3.12) in the space $\dot{X}_{\mu, h}^{1 / 2}$.

Proof. From Salo [18, Lemma 4.1], we know there exists $\sigma$ which satisfies $1 / 2>\sigma>$ $\sigma_{0}, \theta=\sigma-\sigma_{0}$, and symbols $a, b, r_{0},\langle x\rangle r_{0} \in S_{\sigma_{0}, \sigma}$ which give intertwining operators that satisfy

$$
\begin{equation*}
\left(Q+2 h W^{\sharp} \cdot(h D+\mu)\right) A=B Q+h^{1+\tau} R_{0} \tag{3.14}
\end{equation*}
$$

where $\tau=\min \{\theta, 1-2 \sigma\}$, and $a=e^{i \phi}$ with

$$
\left|\partial_{x}^{\alpha} \partial_{\xi}^{\beta} \phi(x, \xi)\right| \leq C_{\alpha, \beta} h^{-\sigma_{0}|\alpha|-\sigma|\beta|}\langle x\rangle^{-1} .
$$

Here, the support of $r_{0}$ is contained in the set $\left\{(x, \xi):|x|<M h^{-\theta}\right\}$. Next we show $A$ and $B$ are bounded and invertible on $X_{\mu, h, \sigma_{0}}^{b}$, for $-1<b<1$ and $\sigma_{0} \in[0,1)$, and when $h$ is small enough, the norms are bounded uniformly in $h$. This follows from Lemma 2.9, since $\frac{1}{a}=e^{-i \phi} \in S_{\sigma_{0}, \sigma}$, then

$$
\mathrm{Op}_{h}(a) \mathrm{Op}_{h}(1 / a)=I+h^{1-\sigma_{0}-\sigma} \mathrm{Op}_{h}(d)
$$

where $d \in S_{\sigma_{0}, \sigma}$. Since $\mathrm{Op}_{h}(d)$ is bounded on $X_{\mu, h, \sigma_{0}}^{b}$, then operator $I+h^{1-\sigma_{0}-\sigma} \mathrm{Op}_{h}(d)$ is invertible on $X_{\mu, h, \sigma_{0}}^{b}$ if $h$ is small enough. So $A$ is invertible with norm of the inverse uniformly bounded in $h$. We can use the same method to prove this property for $B$.

Using the intertwining operators, equation (3.12 becomes

$$
\begin{equation*}
\left(B Q+h^{1+\tau} R_{0}\right) A^{-1} v=g \tag{3.15}
\end{equation*}
$$

Furthermore, (3.12) is equivalent to

$$
\begin{equation*}
\left(I+h^{1+\tau} A Q^{-1} B^{-1} R_{0} A^{-1}\right) v=A Q^{-1} B^{-1} g \tag{3.16}
\end{equation*}
$$

Now we prove that

$$
\begin{equation*}
\left\|A Q^{-1}\right\|_{\dot{X}_{\mu, h}^{-1 / 2}, \dot{X}_{\mu, h}^{1 / 2}} \leq C h^{-\sigma_{0} / 2} \tag{3.17}
\end{equation*}
$$

which is equivalent to

$$
\begin{equation*}
\|A\|_{\dot{X}_{\mu, h}^{1 / 2}, \dot{X}_{\mu, h}^{1 / 2}} \leq C h^{-\sigma_{0} / 2} \tag{3.18}
\end{equation*}
$$

By the estimate of symbol $a$ after (3.14), we have

$$
(a-1)\langle x\rangle \in S_{\sigma_{0}, \sigma}
$$

Thus we can apply Lemma 3.2 to $A$ and obtain (3.18).
Thus, from estimate (3.17), we obtain

$$
\begin{align*}
& \left\|h^{1+\tau} A Q^{-1} B^{-1} R_{0} A^{-1} v\right\|_{\dot{X}_{\mu, h}^{1 / 2}}  \tag{3.19}\\
& \quad \leq C h^{1+\tau-\sigma_{0} / 2}\left\|B^{-1} R_{0} A^{-1} v\right\|_{\dot{X}_{\mu, h}^{-1 / 2}} .
\end{align*}
$$

By estimate (2.10) in Proposition 2.2 and part (3) of Proposition 2.10 we have

$$
\begin{align*}
& \leq C h^{1+\tau-\sigma_{0} / 2}\left\|\langle x\rangle^{1 / 2+\delta} B^{-1} R_{0} A^{-1} v\right\|_{X_{\mu, h}^{-1 / 2}} \\
& \leq C h^{1 / 2+\tau-\sigma_{0} / 2}\left\|\langle x\rangle^{1 / 2+\delta} B^{-1} R_{0} A^{-1} v\right\|_{L^{2}}  \tag{3.20}\\
& \leq C h^{1 / 2+\tau-\sigma_{0} / 2}\left\|B^{-1}\langle x\rangle^{1 / 2+\delta} R_{0}\langle x\rangle^{1 / 2+\delta}\langle x\rangle^{-1 / 2-\delta} A^{-1} v\right\|_{L^{2}}
\end{align*}
$$

By the boundedness of $B^{-1}$ in $L^{2}$ and part (3) of Proposition 2.10, we obtain

$$
\begin{equation*}
\leq C h^{1 / 2+\tau-\sigma_{0} / 2}\left\|\langle x\rangle^{1 / 2+\delta} R_{0}\langle x\rangle^{1 / 2+\delta} A^{-1}\langle x\rangle^{-1 / 2-\delta} v\right\|_{L^{2}} \tag{3.21}
\end{equation*}
$$

Using that the support of $r_{0}(\cdot, \cdot, h)$ is contained in $\left\{(x, \xi):|x|<M h^{-\theta}\right\}$ and that $\langle x\rangle r_{0} \in S_{\sigma_{0}, \sigma}^{0}$, it follows that

$$
\begin{equation*}
\left\|\langle x\rangle^{1 / 2+\delta} R_{0}\langle x\rangle^{1 / 2+\delta}\right\|_{L^{2} \rightarrow L^{2}} \leq C h^{-2 \theta \delta} . \tag{3.22}
\end{equation*}
$$

Finally, the estimate $(3.22)$, the boundedness of $A^{-1}$ followed by the estimate (2.11), gives the bound

$$
\begin{align*}
& \leq C h^{1 / 2+\tau-\sigma_{0} / 2-2 \theta \delta}\left\|\langle x\rangle^{-1 / 2-\delta} v\right\|_{L^{2}} \\
& \leq C h^{\tau-\sigma_{0} / 2-2 \theta \delta}\|v\|_{\dot{X}_{\mu, h}^{1 / 2}} . \tag{3.23}
\end{align*}
$$

Since $\delta$ could be any small positive number, we can choose suitable $\tau, \sigma_{0}$ to make the power of $h$ positive in (3.23) with the requirement $\tau=\min \{\theta, 1-2 \sigma\}, \sigma=\sigma_{0}+\theta$. For example, if we let $\sigma_{0}=\theta$ be some small number and $\delta<1 / 4$, then $\sigma=2 \sigma_{0}$, $\tau=\sigma_{0}$, thus $\tau-\sigma_{0} / 2-2 \theta \delta>\sigma_{0}-\sigma_{0} / 2-\sigma_{0} / 2=0$.

Then by the contraction mapping theorem, there exists a solution $v$ for equation (3.12), and the solution satisfies

$$
\begin{align*}
\|v\|_{\dot{X}_{\mu, h}^{1 / 2}} & \leq C\left\|A Q^{-1} B^{-1} g\right\|_{\dot{X}_{\mu, h}^{1 / 2}} \\
& \leq C h^{-\sigma_{0} / 2}\left\|B^{-1} g\right\|_{\dot{X}_{\mu, h}^{-1 / 2}} \\
& \leq C h^{-\sigma_{0} / 2}\left\|B^{-1}\langle x\rangle^{1 / 2+\delta} g\right\|_{X_{\mu, h}^{-1 / 2}}  \tag{3.24}\\
& \leq C h^{-\sigma_{0}}\left\|\langle x\rangle^{1 / 2+\delta} g\right\|_{X_{\mu, h}^{-1 / 2}}
\end{align*}
$$

For the uniqueness, suppose we have two solutions of (3.12), $v_{1}$ and $v_{2}$, which lie in $\dot{X}_{\mu, h}^{1 / 2}$. Using the intertwining operators, we obtain

$$
\begin{equation*}
\left(B Q+h^{1+\tau} R_{0}\right) A^{-1}\left(v_{1}-v_{2}\right)=0 \tag{3.25}
\end{equation*}
$$

Since Lemma 3.2 gives that $A^{-1}$ is invertible on $\dot{X}_{\mu, h}^{1 / 2}$, it suffices to show that $\tilde{v}=$ $A^{-1}\left(v_{1}-v_{2}\right)$ which satisfies

$$
\left(B Q+h^{1+\tau} R_{0}\right) \tilde{v}=0, \quad \tilde{v} \in X_{\mu, h}^{1 / 2}
$$

is zero. We will show that $\langle\cdot\rangle^{1 / 2+\delta} R_{0} \tilde{v}$ is in $L^{2}$ and then since the equation $Q \tilde{v}=f$ has a unique solution in $\dot{X}_{\mu, h}^{1 / 2}$ when $f$ is in the weighted $L^{2}$-space, $L_{1 / 2+\delta}^{2}$, we have

$$
\begin{equation*}
\left(I+h^{1+\tau} Q^{-1} B^{-1} R_{0}\right) \tilde{v}=0 . \tag{3.26}
\end{equation*}
$$

Thus by (3.22) and (2.11),

$$
\begin{aligned}
\left\|\langle x\rangle^{1 / 2+\delta} R_{0} \tilde{v}\right\|_{L^{2}} & =\left\|\langle x\rangle^{1 / 2+\delta} R_{0}\langle x\rangle^{1 / 2+\delta}\langle x\rangle^{-1 / 2-\delta} \tilde{v}\right\|_{L^{2}} \\
& \leq C h^{-2 \theta \delta}\left\|\langle x\rangle^{-1 / 2-\delta} \tilde{v}\right\|_{L^{2}} \\
& \leq C h^{-2 \theta \delta-1 / 2}\|\tilde{v}\|_{\dot{X}_{\mu, h}^{1 / 2}} .
\end{aligned}
$$

Finally, we use that the operator $B^{-1}$ is bounded on the weighted $L^{2}$-space $L_{1 / 2+\delta}^{2}$, 2.10, and that $Q^{-1}: \dot{X}_{\mu, h}^{-1 / 2} \rightarrow \dot{X}_{\mu, h}^{1 / 2}$ and we have

$$
h^{1+\tau}\left\|Q^{-1} B^{-1} R_{0} \tilde{v}\right\|_{\dot{X}_{\mu, h}^{1 / 2}} \leq C h^{\tau-2 \theta \delta}\|\tilde{v}\|_{\dot{X}_{\mu, h}^{1 / 2} .} .
$$

As $\delta>0$, may be arbitrarily small, it follows that the equation (3.26) has only the solution $\tilde{v}=0$.

Our next theorem gives the existence of CGO solutions to $H_{W, p} u=0$.
Theorem 3.4. Suppose $W \in C^{\lambda}, \lambda>1 / 2, p=p_{1}+\operatorname{div}_{2}, p_{1}, p_{2} \in C^{\tilde{\lambda}}, \tilde{\lambda}>0$; Then for each $\zeta$ large enough, the equation $H_{W, p} u=0$ has a unique CGO solution $u=e^{i x \cdot \zeta}(1+v)$ with $v$ in $\dot{X}_{\mu, h}^{1 / 2}$, and

$$
\|v\|_{\dot{X}_{\mu, h}^{1 / 2}} \leq C h^{2-\sigma_{0}}\left\|\langle x\rangle^{1 / 2+\delta} f\right\|_{X_{\mu, h}^{-1 / 2}} \leq h^{\tilde{\lambda}-\sigma_{0}}
$$

for some $\sigma_{0}>0$ and for any $\delta>0$, and $f$ is defined (3.1.
Proof. We know that $H_{W, p} u=0$ is equivalent to the following equation for $v$

$$
\begin{equation*}
\left(Q+2 h W^{\sharp} \cdot(h D+\mu)+2 h W^{b} \cdot(h D+\mu)+h^{2} G\right) v=h^{2} f . \tag{3.27}
\end{equation*}
$$

Furthermore, as in Lemma 3.3 equation (3.16), the equation (3.27) is equivalent to

$$
\begin{align*}
& \left(I+h^{1+\tau} A Q^{-1} B^{-1} R_{0} A^{-1}+2 h A Q^{-1} B^{-1} W^{b}(h D+\mu)+h^{2} A Q^{-1} B^{-1} G\right) v  \tag{3.28}\\
& =-h^{2} A Q^{-1} B^{-1} f .
\end{align*}
$$

In the argument leading up to 3.23 in Lemma 3.3, we have proved

$$
\begin{align*}
&\left\|h^{1+\tau} A Q^{-1} B^{-1} R_{0} A^{-1} v\right\|_{\dot{X}_{\mu, h}^{1 / 2}} \\
& \leq C h^{\tau-\sigma_{0} / 2-2 \theta \delta}\|v\|_{\dot{X}_{\mu, h}^{1 / 2}} \tag{3.29}
\end{align*}
$$

with suitable $\tau, \sigma_{0}$ to make the power of $h$ positive.
Now for the second remainder term, similarly, by equation (3.17) and the definition of the space $\dot{X}_{\mu, h}^{-1 / 2}$ and $X_{\mu, h}^{1 / 2}$, we have

$$
\begin{align*}
\| h A Q^{-1} B^{-1} & W^{\mathrm{b}}(h D+\mu) v \|_{\dot{X}_{\mu, h}^{1 / 2}} \\
& \leq C h^{1-\sigma_{0} / 2}\left\|B^{-1} W^{\mathrm{b}}(h D+\mu) v\right\|_{\dot{X}_{\mu, h}^{-1 / 2}}  \tag{3.30}\\
& \leq C h^{1-\sigma_{0} / 2}\left\|\langle x\rangle^{1 / 2+\delta} B^{-1} W^{\mathrm{b}}(h D+\mu) v\right\|_{X_{\mu, h}^{-1 / 2}},
\end{align*}
$$

by part 3 of Proposition 2.10 and $\left\|W^{b}\right\|_{\infty} \leq C h^{\sigma_{0} \lambda}, W^{b}$ is supported in a ball $B(0, M)$, we obtain (3.30) is less then

$$
\begin{align*}
& C h^{1 / 2-\sigma_{0} / 2}\left\|B^{-1}\langle x\rangle^{1 / 2+\delta} W^{\mathrm{b}}(h D+\mu) v\right\|_{L^{2}} \\
& \quad \leq C h^{1 / 2-\sigma_{0} / 2}\left\|W^{\mathrm{b}}(h D+\mu) v\right\|_{L^{2}} \\
& \quad \leq C h^{1 / 2-\sigma_{0} / 2}\left\|W^{\mathrm{b}}\right\|_{\infty}\left\|(h D+\mu)\langle x\rangle^{-1 / 2-\delta} v\right\|_{L^{2}}  \tag{3.31}\\
& \quad \leq C h^{1 / 2-\sigma_{0} / 2+\sigma_{0} \lambda}\left\|(h D+\mu)\langle x\rangle^{-1 / 2-\delta} v\right\|_{L^{2}},
\end{align*}
$$

because $|h \xi+\mu| \leq h^{-1 / 2}(h+|q(h \xi)|)^{1 / 2}$, we obtain

$$
\begin{align*}
& \leq C h^{-\sigma_{0} / 2+\sigma_{0} \lambda}\left\|\langle x\rangle^{-1 / 2-\delta} v\right\|_{X_{\mu, h}^{1 / 2}} \\
& \leq C h^{\sigma_{0}(\lambda-1 / 2)}\|v\|_{\dot{X}_{\mu, h}^{1 / 2} .} \tag{3.32}
\end{align*}
$$

Finally, for the last term, by same argument of $A Q^{-1} B^{-1}$ and $\|G\|_{\dot{X}_{\mu, h}^{1 / 2}, \dot{X}_{\mu, h}^{-1 / 2}}<$ $h^{-2+\tilde{\lambda}}$ from Haberman and Tataru [10, Theorem 2.1], then

$$
\begin{align*}
h^{2}\left\|A Q^{-1} B^{-1} G v\right\|_{\dot{X}_{\mu, h}^{1 / 2}} & \leq h^{2-\sigma_{0} / 2}\|G v\|_{\dot{X}_{\mu, h}^{-1 / 2}} \\
& \leq h^{\tilde{\lambda}-\sigma_{0} / 2}\|v\|_{\dot{X}_{\mu, h}^{1 / 2}} \tag{3.33}
\end{align*}
$$

thus, if we choose $\sigma_{0}$ small enough, the power of $h$ will be positive.
Above all, by contraction mapping theorem, there exists a solution $v$ for equation (3.1). And $v$ satisfies

$$
\begin{align*}
\|v\|_{\dot{X}_{\mu, h}^{1 / 2}} & \leq C h^{2}\left\|A Q^{-1} B^{-1} f\right\|_{\dot{X}_{\mu, h}^{1 / 2}} \\
& \leq C h^{2-\sigma_{0} / 2}\left\|B^{-1} f\right\|_{\dot{X}_{\mu, h}^{-1 / 2}} \\
& \leq C h^{2-\sigma_{0} / 2}\left\|B^{-1}\langle x\rangle^{1 / 2+\delta} f\right\|_{X_{\mu, h}^{-1 / 2}}  \tag{3.34}\\
& \leq C h^{2-\sigma_{0}}\left\|\langle x\rangle^{1 / 2+\delta} f\right\|_{X_{\mu, h}^{-1 / 2}}
\end{align*}
$$

Since $f=-(2 W \cdot \zeta+G)$, we have

$$
\begin{align*}
\|v\|_{\dot{X}_{\mu, h}^{1 / 2}} & \leq C h^{2-\sigma_{0}}\left\|\langle x\rangle^{1 / 2+\delta}(2 W \cdot \zeta+G)\right\|_{X_{\mu, h}^{-1 / 2}} \\
& \leq C h^{2-\sigma_{0}}\left(\left\|\langle x\rangle^{1 / 2+\delta} W \cdot \zeta\right\|_{X_{\mu, h}^{-1 / 2}}+\left\|\langle x\rangle^{1 / 2+\delta} G\right\|_{X_{\mu, h}^{-1 / 2}}\right)  \tag{3.35}\\
& \leq C h^{2-\sigma_{0}}\left(h^{-3 / 2}+h^{-2+\tilde{\lambda}}\right) \\
& \leq C h^{\tilde{\lambda}-\sigma_{0}}
\end{align*}
$$

The uniqueness of the solution is similar to our earlier proofs.
Note that we can pick $\sigma_{0}$ arbitrarily small.

## Chapter 4 Existence of solutions in semiclassical spaces

Now we still consider equation

$$
\begin{equation*}
H_{W, p}\left(e^{i x \cdot \zeta}(1+v)\right)=0 \tag{4.1}
\end{equation*}
$$

or equivalently

$$
\begin{equation*}
\left(Q(h D)+2 h W(h D+\mu)+h^{2} G\right) v=h^{2} f \tag{4.2}
\end{equation*}
$$

with $G=W \cdot W+D \cdot W+p, p=p_{1}+\operatorname{div} p_{2}=p^{\sharp}+p^{b}$ with $p^{\sharp}$ smooth and $p^{b} \in \operatorname{div} C^{\tilde{\lambda}}\left(\mathbf{R}^{n} ; \mathbf{C}^{n}\right), f=-(2 W \cdot \zeta+G)$.

But instead of working in the $X$-spaces, this time we discuss the solution in semiclassical Sobolev spaces $H_{\rho, h}^{t}$, which are defined by $\left\{f:\|f\|_{H_{\rho, h}^{t}}=\left\|\langle h D\rangle^{t}\langle x\rangle^{\rho} f\right\|_{2}<\right.$ $\infty\}$, for any real number $t, \rho$.

Like Lemma 3.3, we also have an estimate for operator $Q(h D)+2 h W^{\sharp} \cdot(h D+\mu)$ here, where $W^{\sharp}$ is as defined in Chapter 3.

Lemma 4.1. Equation

$$
\begin{equation*}
\left(Q(h D)+2 h W^{\sharp} \cdot(h D+\mu)\right) v=g \tag{4.3}
\end{equation*}
$$

has a unique solution $v \in H_{\rho, h}^{1}$, for $\rho \in(-1,0)$. Furthermore, we have

$$
\begin{equation*}
\|v\|_{H_{\rho, h}^{1}} \leq C h^{-1}\|g\|_{H_{\rho+1}^{-1}} \leq C h^{-1}\|g\|_{L_{\rho+1}^{2}} \tag{4.4}
\end{equation*}
$$

Proof. By (3.14) and following the argument in the previous section, equation (4.3) becomes

$$
\begin{equation*}
\left(I+h^{1+\tau} A Q^{-1} B^{-1} R_{0} A^{-1}\right) v=A Q^{-1} B^{-1} g \tag{4.5}
\end{equation*}
$$

where the symbols of $A, B, R_{0}$ satisfy $a, b, r_{0},\langle x\rangle r_{0} \in S_{\sigma_{0}, \sigma}$.
From Salo [18, Proposition 2.2], we can see that the operators in $O p_{h}\left(S_{\sigma_{0}, \sigma}\right)$ are bounded in the space $H_{\rho, h}^{t}$ for $-1<\rho<0$ and any real number $t$. And in Salo [18, Propositon 4.1], consider the case when $s=1$, the operator norm of the operator $Q^{-1}$ from $L_{\rho+1}^{2}$ to $H_{\rho, h}^{1}$ is bounded by $C h^{-1}$.

Thus, for the second term in the equation (4.5), and we have

$$
\begin{align*}
\| h^{1+\tau} A Q^{-1} & B^{-1} R_{0} A^{-1} v \|_{H_{\rho, h}^{1}} \\
& \leq C h^{1+\tau}\left\|Q^{-1} B^{-1} R_{0} A^{-1} v\right\|_{H_{\rho, h}^{1}},  \tag{4.6}\\
& \leq C h^{\tau}\left\|B^{-1} R_{0} A^{-1} v\right\|_{L_{\rho+1}^{2}}
\end{align*}
$$

since $\langle x\rangle r_{0} \in S_{\sigma_{0}, \sigma}^{0}$, then we obtain

$$
\begin{align*}
& \leq C h^{\tau}\left\|R_{0} A^{-1} v\right\|_{L_{\rho+1}^{2}} \\
& \leq C h^{\tau}\left\|A^{-1} v\right\|_{L_{\rho}^{2}} \\
& \leq C h^{\tau}\|v\|_{L_{\rho}^{2}}  \tag{4.7}\\
& \leq C h^{\tau}\|v\|_{H_{\rho, h}^{1}}
\end{align*}
$$

So by the contraction mapping theorem, there exists solution to (4.5). Since $\|\cdot\|_{H_{\rho+1}^{-1}} \leq C\|\cdot\|_{L_{\rho+1}^{2}}$, thus the operator norm of $Q^{-1}$ is a map from $H_{\rho+1}^{-1}$ to $H_{\rho+1}^{1}$ is bounded by $\mathrm{Ch}^{-1}$, then the solution satisfies the following estimate

$$
\begin{align*}
\|v\|_{H_{\rho, h}^{1}} & \leq C\left\|A Q^{-1} B^{-1} g\right\|_{H_{\rho, h}^{1}} \\
& \leq C\left\|Q^{-1} B^{-1} g\right\|_{H_{\rho+1}^{-1}} \\
& \leq C h^{-1}\left\|B^{-1} g\right\|_{H_{\rho+1}^{-1}}  \tag{4.8}\\
& \leq C h^{-1}\|g\|_{H_{\rho+1}^{-1}} .
\end{align*}
$$

For the uniqueness, suppose we have two solutions $v_{1}, v_{2} \in H_{\rho, h}^{1}$ of equation 4.3). Apply the intertwining operators, we have then

$$
\begin{equation*}
\left(B Q+h^{1+\tau} R_{0}\right) A^{-1}\left(v_{1}-v_{2}\right)=0 \tag{4.9}
\end{equation*}
$$

If we following the argument in Lemma 3.3 , we can see $A^{-1}$ is invertible on $H_{\rho, h}^{1}$. Let $\tilde{v}=A^{-1}\left(v_{1}-v_{2}\right)$, if we can prove $\tilde{v}=0$, the uniqueness is obtained, where $\tilde{v}$ satisfies

$$
\begin{equation*}
\left(B Q+h^{1+\tau} R_{0}\right) \tilde{v}=0 . \quad \tilde{v} \in H_{\rho, h}^{1} \tag{4.10}
\end{equation*}
$$

Since the operator norm of $Q^{-1}$ from $L_{\rho+1}^{2}$ to $H_{\rho, h}^{1}$ is bounded by $C h^{-1}$ and $\left\|R_{0} \tilde{v}\right\|_{L_{\rho+1}^{2}} \leq$ $\left\|\langle x\rangle^{-1} R_{0} \tilde{v}\right\|_{L_{\rho}^{2}} \leq\|\tilde{v}\|_{L_{\rho}^{2}} \leq\|\tilde{v}\|_{H_{\rho, h}^{1}}$, thus we have

$$
\begin{equation*}
\left(I+h^{1+\tau} Q^{-1} B^{-1} R_{0}\right) \tilde{v}=0 \tag{4.11}
\end{equation*}
$$

Because

$$
\begin{align*}
\left\|h^{1+\tau} Q^{-1} B^{-1} R_{0} \tilde{v}\right\|_{H_{\rho, h}^{1}} & \leq h^{\tau}\left\|B^{-1} R_{0} \tilde{v}\right\|_{L_{\rho+1}^{2}} \\
& \leq h^{\tau}\left\|R_{0} \tilde{v}\right\|_{L_{\rho+1}^{2}}  \tag{4.12}\\
& \leq h^{\tau}\|\tilde{v}\|_{H_{\rho, h}^{1}} .
\end{align*}
$$

Then equation (4.11) has only solution $\tilde{v}=0$.
Theorem 4.2. Equation (4.2) has a unique solution in $H_{\rho, h}^{1}$ for $-1<\rho<0$ and satisfies

$$
\|v\|_{H_{\rho, h}^{1}} \leq C h\|f\|_{H_{\rho+1}^{-1}}
$$

Proof. Let $K=\left(Q+2 W^{\sharp} \cdot(h D+\mu)\right)^{-1}$, equation 4.2 becomes

$$
\begin{equation*}
\left(I+2 h K W^{b} \cdot(h D+\mu)+h^{2} K G\right) v=h^{2} K f \tag{4.13}
\end{equation*}
$$

or in more detail

$$
\begin{gather*}
\left(I+2 h K W^{b} \cdot(h D+\mu)+h^{2} K\left(W \cdot W+D \cdot W^{\sharp}+p^{\sharp}\right)+h^{2} K\left(D \cdot W^{b}+p^{b}\right)\right) v \\
=-h^{2} A Q^{-1} B^{-1} f . \tag{4.14}
\end{gather*}
$$

where $a, b, r_{0},\langle x\rangle r_{0} \in S_{\sigma_{0}, \sigma}^{0}$. Now we consider the terms on the left side of equation (4.14). By Lemma 4.1, we have

$$
\begin{align*}
\| h K W^{b}(h D+ & \mu) v \|_{H_{\rho, h}^{1}} \\
& \leq C\left\|W^{b}(h D+\mu) v\right\|_{L_{\rho+1}^{2}}  \tag{4.15}\\
& \leq C\left\|W^{b}\right\|_{\infty}\|(h D+\mu) v\|_{L_{\rho}^{2}} \\
& \leq C\left\|W^{b}\right\|_{\infty}\|v\|_{H_{\rho, h}^{1}}
\end{align*}
$$

and since $W \cdot W+D \cdot W^{\sharp}+p^{\sharp}$ has compact support, we obtain

$$
\begin{align*}
\| h^{2} K\left(W \cdot W+D \cdot W^{\sharp}+\right. & \left.p^{\sharp}\right) v \|_{H_{\rho, h}^{1}} \\
& \leq C h\left\|\left(W \cdot W+D \cdot W^{\sharp}+p^{\sharp}\right) v\right\|_{L_{\rho+1}^{2}} \\
& \leq C h\left\|\left(W \cdot W+D \cdot W^{\sharp}+p^{\sharp}\right)\right\|_{\infty}\|v\|_{L_{\rho}^{2}}  \tag{4.16}\\
& \leq C h\left\|\left(W \cdot W+D \cdot W^{\sharp}+p^{\sharp}\right)\right\|_{\infty}\|v\|_{H_{\rho, h}^{1}} \\
& \leq C h\|v\|_{H_{\rho, h}^{1}} .
\end{align*}
$$

Since $\|\phi(x)\|_{H_{\rho+1}^{-1}} \leq h^{-1}\|\phi(x)\|_{H_{\rho, h}^{1}}$ for any compact supported function $\phi(x)$, thus

$$
\begin{align*}
\| h^{2} K\left(D \cdot W^{b}+\right. & \left.p^{b}\right) v \|_{H_{\rho, h}^{1}} \\
& \leq C h\left\|\left(D \cdot W^{b}+p^{b}\right) v\right\|_{H_{\rho+1}^{-1}}  \tag{4.17}\\
& \leq C\left\|D \cdot W^{b}+(\operatorname{div})^{-1} p^{b}\right\|_{\infty}\|v\|_{H_{\rho, h}^{1}}
\end{align*}
$$

Above all and by the contraction mapping theorem, there exists a solution to (4.14), and the solution $v$ satisfies that

$$
\begin{align*}
\|v\|_{H_{\rho, h}^{1}} & \leq C h^{2}\left\|A Q^{-1} B^{-1} f\right\|_{H_{\rho, h}^{1}}  \tag{4.18}\\
& \leq C h\|f\|_{H_{\rho+1}^{-1}}
\end{align*}
$$

For the uniqueness, it follows similar argument in Lemma 4.1.

## Chapter 5 From $\Lambda_{W, p}$ to boundary value of CGO solution

In this section, we show the boundary value of solution of $H_{W, p} u=0$ is determined by the DN map. We begin by giving several equivalent characterizations of the CGO solution. This argument follows closely the work of Nachman [14] and Salo [18.

Proposition 5.1. Let $\Omega \subset \mathbf{R}^{n}$ be a bounded domain with $C^{1,1}$ boundary. Suppose $W \in C^{\lambda}, p=p_{1}+\operatorname{div} p_{2}$, with $p_{1}, p_{2} \in C^{\tilde{\lambda}}, \lambda, \tilde{\lambda} \in(1 / 2,1)$. If 0 is not a Dirichlet eigenvalue of $H_{W, p}$ in $\Omega$, let $\zeta \in \mathbf{C}^{n}$ with $\zeta \cdot \zeta=0$, then the following four problems are equivalent.

$$
\begin{align*}
& (D E)\left\{\begin{array}{l}
H_{W, p} u=0 \text { in } \mathbf{R}^{n} \\
u=e^{i \zeta \cdot x}(1+\omega) \text { with } \omega \in \dot{X}_{\mu, h}^{1 / 2}
\end{array}\right. \\
& (I E)\left\{\begin{array}{l}
u+G_{\zeta} \circ(2 W \cdot D u+(W \cdot W+D \cdot W+p) u)=e^{i \zeta \cdot x} \in \mathbf{R}^{n} \\
u \in H_{l o c}^{1}\left(\mathbf{R}^{n}\right)
\end{array}\right. \\
& (E P)\left\{\begin{array}{l}
\text { i) } \Delta u=0 \in \Omega^{\prime} \\
\text { ii) } u \in H^{1}\left(\Omega_{R}^{\prime}\right) \text { for any } R>R_{0} \\
\text { iii) } u \text { satisfies } \sqrt{5.2}) \text { for a.e. } x \in \mathbf{R}^{n} \\
\text { iv) } \frac{\partial u}{\partial \nu^{+}}=\Lambda_{W, p}\left(u_{+}\right) \text {on } \partial \Omega
\end{array}\right.  \tag{5.1}\\
& (B E)\left\{\begin{array}{l}
\left(\frac{1}{2} I+S_{\zeta} \Lambda_{W, p}-B_{\zeta}\right) f=e^{i \zeta \cdot x} \text { on } \partial \Omega \\
f \in H^{1 / 2}(\partial \Omega)
\end{array}\right.
\end{align*}
$$

We first show the differential equation (DE) and the integral equation (IE) are equivalent. Now define the Green function $G_{\zeta}$ by

$$
G_{\zeta}=e^{i \zeta \cdot x} g_{\zeta} e^{-i \zeta \cdot x}
$$

where $g_{\zeta}$ is the fundamental solution to $\Delta_{\zeta}=e^{-i \zeta \cdot x} \Delta e^{i \zeta \cdot x}$ such that $\Delta_{\zeta}^{-1} f=g_{\zeta} * f$ for $f$ in the Schwartz class. Then

$$
\Delta G_{\zeta}=\zeta^{2} G_{\zeta}+2 e^{i \zeta \cdot x} \zeta \cdot D g_{\zeta}+e^{i \zeta \cdot x} \Delta g_{\zeta}=e^{i \zeta \cdot x} \Delta_{\zeta} g_{\zeta}=\delta_{0}
$$

where $\delta_{0}$ is the Dirac measure at 0 . So we can let $G_{\zeta}=G_{0}+H_{\zeta}$ where $G_{0}(x)=c_{n}|x|^{2-n}$ is the fundamental solution of $\Delta, c_{n}=\frac{1}{n(n-2) \alpha(n)}, \alpha(n)$ is the volume of unit ball in $\mathbf{R}^{n}$ and $H_{\zeta}$ is a global harmonic function.

Lemma 5.2. Suppose we have the same conditions as Proposition 5.1. Then $u$ is a solution of ( $D E$ ) if and only if $u$ is a solution of (IE). Also, a solution of (DE) is unique if and only if $u$ is a solution of (IE) is unique.

Proof. Let $u=e^{i \zeta \cdot x}(1+\omega)$ solve (DE) where $\omega=\Delta_{\zeta}^{-1} f$ with $f \in \dot{X}_{\mu, h}^{1 / 2}$. Substitute $u$ into $H_{W, p} u=0$, we have

$$
\left(\Delta_{\zeta}+2 W \cdot D_{\zeta}+(W \cdot W+D \cdot W+p)\right)\left(1+\Delta_{\zeta}^{-1} f\right)=0
$$

Since $\Delta_{\zeta}\left(1+\Delta_{\zeta}^{-1} f\right)=f$, apply $\Delta_{\zeta}^{-1}$ on both sides; then

$$
\omega+\Delta_{\zeta}^{-1}\left(2 W \cdot D_{\zeta}(1+\omega)+(W \cdot W+D \cdot W+p)(1+\omega)\right)=0
$$

Next we add one on both sides and multiply $e^{i \zeta \cdot x}$, and we obtain (IE).
If $u$ solves (IE), write $u=e^{i \zeta \cdot x} u_{0}$; then $u_{0}$ solves

$$
u_{0}+\Delta_{\zeta}^{-1}\left(2 W \cdot D_{\zeta} u_{0}+(W \cdot W+D \cdot W+p) u_{0}\right)=1
$$

If we apply $\Delta_{\zeta}$ to both sides, we have $H_{W, p} u=0$.
The uniqueness part is obtained by noting that if $u_{1}$ and $u_{2}$ solve (DE) then $u_{1}$ and $u_{2}$ solve (IE), and vice versa.

Now, we show that (IE) and the exterior problem (EP) are equivalent. We use the notation $\Omega^{\prime}=\mathbf{R}^{n} \backslash \bar{\Omega}$ and $\Omega_{R}^{\prime}=B(0, R) \backslash \bar{\Omega}$, where $R>R_{0}$ and $\bar{\Omega} \subseteq B\left(0, R_{0}\right)$. Let $u_{+}$(resp. $u_{-}$) for the restriction of $u$ to $\partial \Omega$ from the exterior (resp. interior), and $\frac{u}{\partial \nu_{+}}$(resp. $\frac{u}{\partial \nu_{-}}$) for the value of $\nabla u \cdot \nu$ on $\partial \Omega$ from the exterior (resp. interior), where $\nu$ is the outward unit normal to $\partial \Omega$. We also write $G_{\zeta}(x, y)=G_{\zeta}(x-y)$.

We want to obtain that a solution of (IE) satisfies the radiation condition

$$
\begin{equation*}
\int_{|y|=R}\left(G_{\zeta}(x, y) \frac{\partial u}{\partial \nu}(y)-u(y) \frac{\partial G_{\zeta}(x, y)}{\partial \nu(y)} d S(y) \rightarrow e^{i \zeta \cdot x}\right. \tag{5.2}
\end{equation*}
$$

for a.e. $x \in \mathbf{R}^{n}$ as $R \rightarrow \infty$. In order to apply Green's identity, we define a smooth approximation of $G_{\zeta}$ by $G_{\zeta}^{\varepsilon}=G_{0}^{\varepsilon}+H_{\zeta}$, where

$$
G_{0}^{\varepsilon}=c_{n}\left(\varepsilon^{2}+|x|^{2}\right)^{\frac{2-n}{2}} .
$$

In fact $\Delta G_{\zeta}^{\varepsilon}(x)=\varepsilon^{-n} \varphi(x / \varepsilon)$ where

$$
\varepsilon(x)=\frac{1}{\alpha(n)}\left(1+|x|^{2}\right)^{-\frac{n+2}{2}}
$$

and $\int \varphi(x) d x=1$. Then $\Delta G_{\zeta}^{\varepsilon}$ is an approximation of the identity.
We need a lemma on regularity properties of solution $H_{W, p} u=0$ and of $\Lambda_{W, p}$.
Lemma 5.3. Under the conditions of Proposition 5.1, the operator $P_{W, p}$, which maps $f \in H^{1 / 2}(\Omega)$ to the solution $u$ of $H_{W, p} u=0$ in $\Omega$ with $\left.u\right|_{\partial \Omega}=f$ is bounded $H^{1 / 2}(\partial \Omega) \rightarrow H^{1}(\Omega)$. Further, we have $\Lambda_{W, p}: H^{1 / 2}(\partial \Omega) \rightarrow H^{-1 / 2}(\partial \Omega)$, and

$$
\Lambda_{W, p} f=\left.\frac{\partial u}{\partial \nu}\right|_{\partial \Omega} .
$$

Proof. The operator $H_{W, p}$, written in nondivergence form, satisfies the assumption of [8]. This gives that $u$ is in $H^{1}(\Omega)$ if $f \in H^{1 / 2}(\partial \Omega)$, and the solution operator $P_{W, p}$ is bounded.

For the $\Lambda_{W, p}$ part, we claim that if $W \in L_{\Omega}^{n}\left(\mathbf{R}^{n} ; \mathbf{C}^{n}\right)$ and $D \cdot W \in L^{n / 2}\left(\mathbf{R}^{n} ; \mathbf{C}\right)$, then for any $v \in W^{1, n /(n-1)}(\Omega)$ one has

$$
\begin{equation*}
\int_{\Omega}(W \cdot D v+(D \cdot W) v) d x=0 \tag{5.3}
\end{equation*}
$$

This statement means that $W \cdot \nu=0$ on $\partial \Omega$, in some weak sense. Then we take $W_{j} \in C_{c}^{\infty}\left(\mathbf{R}^{n} ; \mathbf{C}^{n}\right)$ to be the convolution approximation of $W$ so that $W_{j} \rightarrow W$ in $L^{n}$ and $D \cdot W_{j} \rightarrow D \cdot W$ in $L^{n / 2}$, and we take an extension of $v$ in $W^{1, n /(n-1)}\left(\mathbf{R}^{n}\right)$. If the supports of $W_{j}$ and $W$ are contained in $B(0, R)$, then

$$
\begin{align*}
\int_{\Omega}(W \cdot D v+(D \cdot W) v d x & =\lim _{j \rightarrow \infty} \int_{B(0, R)}\left(W_{j} \cdot D v+\left(D \cdot W_{j}\right) v\right) d x  \tag{5.4}\\
& =\lim _{j \rightarrow \infty} \frac{1}{i} \int_{\partial B(0, R)}\left(W_{j} \cdot \nu\right) v d S=0 .
\end{align*}
$$

Let $f, g \in H^{1 / 2}(\Omega)$ and $u_{f}=P_{W, p} f$ and $e_{g} \in H^{1}(\Omega)$ with $\left.e_{g}\right|_{\partial \Omega}=g$. Use the definition of $\Lambda_{W, p}$ to obtain

$$
\begin{equation*}
\left\langle\left.\frac{\partial u_{f}}{\partial \nu}\right|_{\partial \Omega)}, g\right\rangle=\int_{\omega}\left(\nabla u_{f} \cdot \nabla e_{g}+\left(2 W \cdot D u_{f}+(W \cdot W+D \cdot W+p) u_{f}\right) e_{g}\right) d x . \tag{5.5}
\end{equation*}
$$

Now $u_{f} e_{g} \in W^{2,1}(\Omega) \subseteq W^{1, n /(n-1)}(\Omega)$. Using (5.3) with $v=u_{f} e_{g}$ and substituting to (5.5) gives $\left.\frac{\partial u_{f}}{\partial \nu}\right|_{\partial \Omega}=\Lambda_{w, p} f$.

Lemma 5.4. Assume the conditions of Proposition (5.1). If $u$ is a solution of (IE), then $\left.u\right|_{\Omega^{\prime}}$ is a solution of (EP). Conversely, if $u$ is a solution of $(E P)$, then there is a unique extension of $u$ to $\mathbf{R}^{n}$ so that $\tilde{u}$ is a solution of (IE). The uniqueness also holds for (IE) and (EP).

Proof. If $u$ solves (IE). By Lemma 5.2, we have $H_{W, p} u=0$ and $u=e^{i x \cdot \zeta}(1+\omega)$ with $\omega \in \dot{X}_{\mu, h}^{1 / 2}$, which gives us (EP) i) - ii). To show iii), for fixed $x$, let $R>|x|$ and $R>R_{0}$, and write

$$
\begin{align*}
\int_{|y|=R} & \left(G_{\zeta}(x, y) \frac{\partial u}{\partial \nu}(y)-u(y) \frac{\partial G_{\zeta}(x, y)}{\partial \nu(y)} d S(y)\right. \\
& =-\int_{B(0, R)}\left(G_{\zeta}^{\varepsilon}(x, y) \frac{\partial u}{\partial \nu}(y)-u(y) \Delta_{y} G_{\zeta}^{\varepsilon}(x, y)\right) d y \\
& =\int_{B(0, R)} u \Delta_{y} G_{\zeta}^{\varepsilon}(x, y) d y+\int_{B(0, R)} G_{\zeta}^{\varepsilon}(x, y)(2 W \cdot D u+(W \cdot W+D \cdot W+p) u) d y \\
& =\left(\Delta G_{\zeta}^{\varepsilon} * u \chi_{B(0, R)}\right)(x)+\left(G_{\zeta}^{\varepsilon}(2 W \cdot D u+(W \cdot W+D \cdot W+p) u)\right)(x) \tag{5.6}
\end{align*}
$$

since $W$ and $p$ supported inside $B(0, R)$, a solution of (IE) is harmonic and in $C^{\infty}\left(\mathbf{R}^{n} \backslash\right.$ $\Omega)$. As $\varepsilon \rightarrow 0$, the first term converges to $u(x)$, a.e. $x$. The second term converges to $\left(G_{\zeta} *(2 W \cdot D u+(W \cdot W+D \cdot w+p) u)\right)(x)$ for a.e. $x$. This gives us (EP) iii). By

Lemma 5.3. since $u \in H^{1}$ and 0 is not a Dirichlet eigenvalue of $H_{W, p}$, we obtain (EP) iv).

If $u$ solve (EP). We use Lemma 5.3 and let $v=P_{W, p} u_{+}$and define $\tilde{u}(x)=u(x)$ for $x \in \Omega^{\prime}$ and $\tilde{u}(x)=v(x)$ for $x \in \Omega$. By (EP) i) we have $H_{W, p} \tilde{u}=0$. The uniqueness part follows from the facts that if $u_{1}$ and $u_{2}$ solve (IE) then $\left.u_{1}\right|_{\Omega^{\prime}}$ and $\left.u_{2}\right|_{\Omega^{\prime}}$ solve (EP), and if $u_{1}$ and $u_{2}$ solve (EP) then $\tilde{u_{1}}$ and $\tilde{u_{2}}$ solve (IE).

To prove (EP) and (BE), we need to use the single layer potential $S_{\zeta}$, double layer potential $D_{\zeta}$, and boundary layer potential $B_{\zeta}$, which are defined by

$$
\begin{align*}
S_{\zeta} f(x) & =\int_{\partial \Omega} G_{\zeta}(x, y) f(y) d S(y) & \left(x \in \mathbf{R}^{n} \backslash \partial \Omega\right) \\
D_{\zeta} f(x) & =\int_{\partial \Omega} \frac{\partial G_{\zeta}(x, y)}{\partial \nu(y)} f(y) d S(y) & \left(x \in \mathbf{R}^{n} \backslash \partial \Omega\right)  \tag{5.7}\\
B_{\zeta} f(x) & =\int_{\partial \Omega} \frac{\partial G_{\zeta}(x, y)}{\partial \nu(y)} f(y) d S(y) & (x \in \partial \Omega)
\end{align*}
$$

Lemma 5.5. Under the condition of Proposition 5.1, if $u$ is a solution of $(E P)$, then $f=\left.u\right|_{\partial \Omega}$ is a solution of $(B E)$. Conversely, if $f$ is a solution of $(B E)$, then

$$
\begin{equation*}
u=e^{i \zeta \cdot x}-S_{\zeta} \Lambda_{W, p} f+D_{\zeta} f \tag{5.8}
\end{equation*}
$$

is a solution of ( $E P$ ), with $u_{+}=f$. Also, solutions of ( $E P$ ) are unique if and only if solutions of ( $B E$ ) are unique.

Proof. Suppose $u$ solves (EP). If we let $f=u_{+}$on $\partial \Omega$. Then $f \in H^{3 / 2}(\partial \Omega)$. If $x \in \Omega^{\prime}$ and $R>|x|$, we have

$$
\begin{align*}
-\int_{\Omega_{R}^{\prime}}\left(G_{\zeta}^{\varepsilon}(x, y)\right. & \left.\Delta u(y)-u(y) \Delta_{y} G_{\zeta}^{\varepsilon}(x, y)\right) d y  \tag{5.9}\\
& =\left(\int_{|y|=R}-\int_{\partial \Omega}\right)\left(G_{\zeta}(x, y) \frac{\partial u}{\partial \nu}(y)-u(y) \frac{\partial G_{\zeta}(x, y)}{\partial \nu(y)} d S(y)\right.
\end{align*}
$$

Let $\varepsilon \rightarrow 0$ and use (EP) i), iii) - iv), we have $f$ solves (BE).
If $f$ satisfies (BE) and define $u$ by (5.8) in $\Omega^{\prime}$. Then $u$ satisfies (EP) i) - iii). We need to show (EP) iv). Let $R \rightarrow \infty$, we can obtain iv).

Since (5.8) gives a correspondence between solutions of (BE) and (EP), this gives the equivalence of uniqueness for the two problems.

Lemma 5.6. Let $\Omega \subset \mathbf{R}^{n}$, $n \geq 3$, be a bounded domain with $C^{1,1}$ boundary. Then the operator $S_{\zeta} \Lambda_{W, p}-B_{\zeta}-\frac{1}{2} I: H^{1 / 2}(\partial \Omega) \rightarrow H^{1 / 2}(\partial \Omega)$ is compact.
Proof. Suppose $f \in H^{1 / 2}(\partial \Omega)$ and $u=P_{W, p} f$. Then for $x \in \Omega$, we have

$$
\begin{align*}
-\int_{\Omega}\left(G_{\zeta}^{\varepsilon}(x, y)\right. & \left.\Delta u(y)-u(y) \Delta_{y} G_{\zeta}^{\varepsilon}(x, y)\right) d y \\
= & \int_{\partial \Omega}\left(G^{\varepsilon} \zeta(x, y) \frac{\partial u}{\partial \nu}(y)-u(y) \frac{\partial G_{\zeta}^{\varepsilon}(x, y)}{\partial \nu(y)} d S(y)\right. \tag{5.10}
\end{align*}
$$

Let $\varepsilon \rightarrow 0$ we obtain

$$
u(x)=\int_{\Omega} G_{\zeta}(x, y)\left(2 W \cdot D+(W \cdot W) u d y+\langle D \cdot W+p, u\rangle=\left(S_{\zeta} \Lambda_{W, p}-D_{\zeta}\right) f(x)\right.
$$

a.e. in $\Omega$. If let $x \rightarrow \partial \Omega$ nontangentially, then we get

$$
\begin{aligned}
& \left(S_{\zeta} \Lambda_{W, p}-B_{\zeta}-\frac{1}{2} I\right) f(x) \\
& \quad=R \int_{\Omega} G_{\zeta}(x, y)(2 W \cdot D+(W \cdot W)) P_{W, p} f(y) d y+\left\langle D \cdot W+p, P_{W, p} f(y)\right\rangle
\end{aligned}
$$

which can be written as

$$
S_{\zeta} \Lambda_{W, p}-B_{\zeta}-\frac{1}{2} I=R G_{\zeta} M P_{W, p}
$$

where $R$ is the trace $H^{1}(\Omega) \rightarrow H^{1 / 2}(\Omega), G_{\zeta}: H^{-1}(\Omega) \rightarrow H^{1}(\Omega)$ is the map that restricts $G_{\zeta} \tilde{u}=e^{i \zeta \cdot x} \Delta_{\zeta}^{-1} e^{-i \zeta \cdot x} \tilde{u}$ to $\Omega$ with $\tilde{u}$ be the extension by zero of $u \in L^{2}(\Omega)$ to $\mathbf{R}^{n}$. The map $M: H^{1} \Omega \rightarrow H^{-1}(\Omega)$ maps $u$ to $2 W \cdot D u+(W \cdot W+D \cdot W+p) u$ since $|\langle p u, v\rangle| \leq\|u\|_{H^{1}}\|v\|_{H^{1}}$. Since $W, p_{1}$ and $p_{2}$ are Hölder continuous, the composition is compact.

Proposition 5.7. Suppose the conditions of Theorem 1.3, then there exists $C=$ $C(n, W, p, \Omega)$ such that for any $|\zeta| \geq C$, each of the four problems (DE), (IE), (EP), (BE) has a unique solution.

Proof. If we can show the problem (DE) has a unique solution, then the other three problems will also have a unique solution. We know problem (DE) has a unique solution from Lemma 3.4.

## Chapter 6 Recovery of curl $W$

Let $u_{\zeta}=e^{i \zeta \cdot x}\left(\omega_{0}+\omega\right)$ be the solution of $H_{W, p} u_{\zeta}=0$ with

$$
\begin{align*}
& \omega_{0}=e^{i \chi_{\zeta} \phi^{\sharp}} \\
& \chi_{\zeta}(x)=\chi\left(x /|\zeta|^{\theta}\right),  \tag{6.1}\\
& \phi^{\sharp}(x)=N_{\mu}^{-1}\left(-\mu \cdot W^{\sharp}\right) .
\end{align*}
$$

where function $\chi \in C_{0}^{\infty}\left(\mathbf{R}^{n}\right), \chi=1$ in $B(0, M / 2), \chi=0$ outside $B(0, M)$, and $\bar{\Omega} \in B(0, M / 2), W^{\sharp}$ is as defined in (3.3); the operator $N_{\mu}=\mu \cdot \nabla$, and when $\mu=\gamma_{1}+i \gamma_{2}$ where $\left|\gamma_{j}\right|=1$ and $\gamma_{1} \cdot \gamma_{2}=0$,

$$
N_{\mu}^{-1} f(x)=\frac{1}{2 \pi} \int_{\mathbf{R}^{2}} \frac{1}{y_{1}+i y_{2}} f\left(x-y_{1} \gamma_{1}-y_{2} \gamma_{2}\right) d y_{1} d y_{2} .
$$

Lemma 6.1. The function $\omega \in H_{\rho, h}^{1}$ for $-1<\rho<0$ and $\|\omega\|_{H_{\rho, h}^{1}}=o(1)$ as $h \rightarrow 0$.
Proof. Substitute $u_{\zeta}$ into $H_{W, p} u_{\zeta}=0$ and we get,

$$
\begin{equation*}
\left(\Delta_{\zeta}+2 W \cdot D_{\zeta}+G\right) \omega=-f \tag{6.2}
\end{equation*}
$$

or equivalently

$$
\begin{equation*}
\left(Q+2 h W \cdot(h D+\mu)+h^{2} G\right) \omega=-h^{2} f \tag{6.3}
\end{equation*}
$$

where $G=W \cdot W 2+D \cdot W+p$ and

$$
\begin{align*}
f= & \left(\Delta_{\zeta}+2 W \cdot D_{\zeta}+G\right) \omega_{0} \\
= & e^{i \chi_{\zeta} \phi^{\sharp}}\left[i \chi_{\zeta} \Delta \phi^{\sharp}+2 i D \chi_{\zeta} \cdot D \phi^{\sharp}+i\left(\Delta \chi_{\zeta}\right) \phi^{\sharp}+\left(\chi_{\zeta} \nabla \phi^{\sharp}+\left(\nabla \chi_{\zeta}\right) \phi^{\sharp}\right)^{2}\right. \\
& +2 \zeta \cdot\left(\nabla \chi_{\zeta}\right) \phi^{\sharp}+2 \zeta \cdot\left(\chi_{\zeta} \nabla \phi^{\sharp}\right)+2 W \cdot\left(\nabla \chi_{\zeta}\right) \phi^{\sharp}+2 W \cdot\left(\chi_{\zeta} \nabla \phi^{\sharp}\right)  \tag{6.4}\\
& \left.+2 W^{\sharp} \cdot \zeta+2 W^{b} \cdot \zeta+G\right],
\end{align*}
$$

Since

$$
\zeta \cdot \nabla \phi^{\sharp}+W^{\sharp} \cdot \zeta=0
$$

and $W^{\sharp}=\chi_{\zeta} W^{\sharp}$, which can cancel two terms of $f$. By Theorem4.2, we have that

$$
\begin{align*}
\|\omega\|_{H_{\rho, h}^{1}} \leq & C h\|f\|_{H_{\rho+1}^{-1}} \\
\leq & C h\left[\left\|\chi_{\zeta} \Delta \phi^{\sharp}\right\|_{L_{\rho+1}^{2}}+\left\|\nabla \chi_{\zeta} \cdot \nabla \phi^{\sharp}\right\|_{L_{\rho+1}^{2}}+\left\|\left(\Delta \chi_{\zeta}\right) \phi^{\sharp}\right\|_{L_{\rho+1}^{2}}\right. \\
& +\left\|\left|\chi_{\zeta} \nabla \phi^{\sharp}\right|^{2}\right\|_{L_{\rho+1}^{2}}+\left\|\left|\left(\nabla \chi_{\zeta}\right) \phi^{\sharp}\right|^{2}\right\|_{L_{\rho+1}^{2}}+|\zeta|^{1-\theta}\left\|\left(\nabla \chi\left(x /|\zeta|^{\theta}\right)\right) \phi^{\sharp}\right\|_{L_{\rho+1}^{2}}  \tag{6.5}\\
& +\left\|W \cdot\left(\nabla \chi_{\zeta}\right) \phi^{\sharp}\right\|_{L_{\rho+1}^{2}}+\left\|W \cdot\left(\chi_{\zeta} \nabla \phi^{\sharp}\right)\right\|_{L_{\rho+1}^{2}}+\left\|W^{b} \cdot \zeta\right\|_{L_{\rho+1}^{2}} \\
& \left.+h\left\|G e^{i \chi_{\zeta} \phi^{\sharp}}\right\|_{H_{\rho, h}^{-1}}\right] .
\end{align*}
$$

The first nine terms in the brackets, are $o\left(h^{-1}\right)$ as in Salo [18]. For the tenth term, since $G=W \cdot W+D \cdot W+p$, we use the same argument as in (4.16) and (4.17), to show it is bounded by $\left\|G e^{i \chi_{\zeta} \phi^{\sharp}}\right\|_{H_{\rho, h}^{-1}} \leq o\left(h^{-\sigma_{0}}\right)$. Combining with the $h$ outside the brackets, we conclude that $\|\omega\|_{H_{\rho, h}^{1}}=o(1)$, as $h \rightarrow 0$.

Now if $\xi \in \mathbf{R}^{n}$ with $|\xi|^{2}$ not a Dirichlet eigenvalue of $\Delta$, then for any $\zeta \in \mathbf{C}^{n}$ satisfies $\zeta \cdot \zeta=0,|\zeta| \geq C>0, \operatorname{Re} \zeta \perp \xi$ and $\operatorname{Im} \zeta \perp \xi$, we define

$$
\begin{align*}
t_{W, p}(\xi, \zeta) & =\left\langle\left(\Lambda_{W, p}-\Lambda_{0,-|\xi|^{2}}\right)\left(u_{\zeta} \mid \partial \Omega\right), e^{-i x \cdot(\xi+\zeta)} \mid \partial \Omega\right\rangle \\
& =\int_{\Omega} e^{-i x \cdot \xi}\left(2(\zeta \cdot W) u_{0}+W \cdot D u_{0}+\left(\xi \cdot W+|\xi|^{2}\right) u_{0}\right) d x+\left\langle e^{-i x \cdot \xi} p, u_{0}\right\rangle \tag{6.6}
\end{align*}
$$

where $u_{\zeta}=e^{i \zeta \cdot x} u_{0}$ satisfies $H_{W, p} u_{\zeta}=0$.
Replacing $u_{0}$ with $\omega+\omega_{0}$, we get

$$
R_{W, p}(\xi, \mu)=\lim _{h \rightarrow 0} h t_{W, p}(\xi, \zeta)=2 \int e^{-i x \cdot \xi} e^{i \phi}(\mu \cdot W) d x
$$

where $\phi=N_{\mu}^{-1}(-\mu \cdot W)$. We use one lemma in Salo's paper [18] [Lemma 6.2]
Lemma 6.2. One has

$$
R_{W, p}(\xi, \mu)=2 \int e^{-i x \cdot \xi}(\mu \cdot W) d x
$$

To recover curl $W$, we need recover $D_{j} W_{k}-D_{k} W_{j}$ for any $j \neq k$. For any $\mu$ with $|\xi|^{2}$ not a Dirichlet eigenvalue of $\Delta$, we let $\mu_{1}=\frac{\xi_{j} e_{k}-\xi_{k} e_{j}}{\left|\xi_{j} e_{k}-\xi_{k} e_{j}\right|}$; then $\mu_{1} \cdot \xi=0$. Find an unit vector $\mu_{2} \in \mathbf{R}^{n}$ with $\mu_{2} \cdot \xi=\mu \cdot \mu_{1}=0$, now if we let $\mu=\mu_{1}+i \mu_{2}$, from Lemma 6.2, we know we can recover $R_{W, q}(\xi, \mu)$ and $R_{W, p}(\xi, \bar{\mu})$, since $R_{W, p}(\xi, \mu)+R_{W, p}(\xi, \bar{\mu})$ determines

$$
\int e^{-i x \cdot \xi}(\mu+\bar{\mu}) \cdot W d x=\left(D_{j} W_{k}-D_{k} W_{j}\right)^{\wedge}(\xi)
$$

Thus we can recover $D_{j} W_{k}-D_{k} W_{j}$ from this Fourier transform.

## Chapter 7 Recovery of $p$

In this chapter, we show how to recover $p$ from the D-N map. This is similar to what we did in the last chapter, but we estimate the CGO solution in Bourgain type spaces. In the previous chapter, we recovered curl $W$; now we need to construct a certain $W$ from curl $W$. Given a vector field $W=\left(W_{1}, \ldots, W_{n}\right)$, we consider the one form $\sum_{i=1}^{n} W_{i} d x_{i}$. Applying the exterior derivative, we obtain

$$
d\left(\sum_{i=1}^{n} W_{i} d x_{i}\right)=\sum_{i<j} \partial_{x_{i}} W_{j} d x_{i} \wedge d x_{j}
$$

Thus, the coefficients on right-hand side are the components of curl $W$. This implies that finding a solution $\widetilde{W}$ of $\operatorname{curl} \widetilde{W}=\operatorname{curl} W$ is equivalent to finding a solution of $d \sum \widetilde{W}_{i} d x_{i}=d \sum_{i=1}^{n} W_{i} d x_{i}$. For the next result, we work with the differential form $\sum_{i=1}^{n} W_{i} d x_{i}$ (which we still denote by $W$ ) and the exterior derivative $d$, rather than the vector $W$ and the operator curl.

Since we assume that $W$ is compactly supported in $\Omega$, we may use a partition of unity to reduce to the case where $W$ is supported in a ball. According to Mitrea, Mitrea and Monniaux [13, Theorem 4.1], we have operators $J_{1}$ and $J_{2}$ so that

$$
\begin{equation*}
W=J_{2}(d W)+d\left(J_{1} W\right) \tag{7.1}
\end{equation*}
$$

Since $d^{2}\left(J_{1} W\right)=0, \widetilde{W}=J_{2}(d W)$ will be give us a solution of $d \widetilde{W}=d W$ as desired. The coefficients of $J_{l}(\omega), l=1,2, \cdots, n$ are defined by

$$
J_{l} u(x)=\int_{\Omega} \int_{1}^{\infty}(t-1)^{n-l} t^{l-1} \varphi(y+t(x-y))(x-y) \vee u(y) d t d y
$$

with $\varphi \in C_{c}^{\infty}(B(0, R))$ with $\int \varphi=1$ and $1 \leq l \leq n$. Let

$$
T_{l, j} f(x)=\int_{\Omega} \int_{1}^{\infty}(t-1)^{n-l} t^{l-1} \varphi(y+t(x-y))\left(x_{j}-y_{j}\right) f(y) d t d y
$$

We can write

$$
T_{l, j}(f(x))=\int_{\Omega} k_{l, j}(x, x-y) f(y) d y
$$

for $1 \leq m \leq n$, where $k_{l, j}$ is the kernel supported in a ball $\{(x, y):|(x, y)|<R\}$ and satisfies

$$
\begin{equation*}
\left|\partial_{x}^{\alpha} \partial_{y}^{\beta} k_{l, j}(x, y)\right| \leq \frac{C}{|y|^{n-1+|\beta|}} \tag{7.2}
\end{equation*}
$$

Finally, to compute $J_{2}(d W)$, we want to consider the map $T_{\ell, j}\left(\partial_{m} u\right)$; and we need to show that this maps into $C^{\lambda}(B)$ when $u$ is in $C_{c}^{\lambda}(B)$.

Lemma 7.1. If $u \in C_{c}^{\infty}(B)$, then

$$
\begin{equation*}
T_{l, j}\left(\partial_{y_{m}} u(x)\right)=-\int_{\mathbf{R}^{n}} \partial_{y_{m}} k_{l, j}(x, x-y)(u(y)-u(x)) d y \tag{7.3}
\end{equation*}
$$

Proof. This is a result of using integration by parts and the fact $u$ is compactly supported.

If $u \in C_{c}^{\lambda}(B)$, we may find $u_{j} \in C_{c}^{\infty}(B)$ with $u_{i} \rightarrow u$ uniformly and $\left\|u_{i}\right\|_{C^{\lambda}} \leq$ $C\|u\|_{C^{\lambda}}$. Then we may write

$$
\begin{align*}
& T_{l, j}\left(\partial_{y_{m}} u\right)(x) \\
& \quad=-\lim _{i \rightarrow \infty} \int \partial_{y_{m}} k_{l, j}(x, x-y)\left(u_{i}(y)-u_{i}(x)\right) d y \tag{7.4}
\end{align*}
$$

Theorem 7.2. If $u \in C_{c}^{\lambda}(\Omega)$, then

$$
\left\|T_{l, j} \partial_{y_{m}} u\right\|_{C^{\lambda}} \leq c\|u\|_{C^{\lambda}}
$$

Proof. We easily have

$$
\left\|T_{l, j} \partial_{y_{m}} u\right\|_{\infty} \leq C\|u\|_{C^{\lambda}}
$$

Then we can write

$$
\begin{align*}
& T_{l, j} \partial_{m} u(y)-T_{l, j} \partial_{m} u(x) \\
& \quad=\int \partial_{z_{m}} k_{l, j}(y, y-z)(u(z)-u(y))-\partial_{z_{m}} k_{l, j}(x, x-z)(u(z)-u(x)) d z \tag{7.5}
\end{align*}
$$

We set $\bar{x}=1 / 2(x+y), s=10|x-y|$, we use number 10 here to make sure we can have a suitable distance between $z$ and $\bar{x}$. Then

$$
\begin{align*}
T_{l, j} \partial_{m}(y)-T_{l, j} \partial_{m} u(x)= & \int_{|\bar{x}-z|<d} \partial_{z_{m}} k_{l, j}(y, y-z)(u(z)-u(y)) d z \\
& -\int_{|\bar{x}-z|<d} \partial_{z_{m}} k_{l, j}(x, x-z)(u(z)-u(x)) d z \\
& +\int_{|\bar{x}-z|>d}\left(\partial_{z_{m}} k_{l, j}(y, y-z)-\partial_{z_{m}} k_{l, j}(x, x-z)\right)  \tag{7.6}\\
& \times(u(z)-u(y)) d z \\
& +(u(x)-u(y)) \int_{|\bar{x}-z|>d} \partial_{z_{m}} k_{l, j}(x, x-z) d z \\
= & I+I I+I I I+I V .
\end{align*}
$$

We have $I+I I \leq C s^{\lambda}$. For $I I I$, we have the estimate

$$
\left.\mid \partial_{z_{m}} k_{l, j}(y, y-z)-\partial_{z_{m}} k_{l, j}(x, x-z)\right) \left\lvert\, \leq C \frac{|x-y|}{|z-\bar{x}|^{n+1}}\right.
$$

for $|z-\bar{x}| \geq s$. Then

$$
\begin{align*}
I I I & \leq C s \int_{s}^{\infty} r^{-n-1} r^{n-1+\lambda} d r \\
& \leq C s \cdot s^{\lambda-1}  \tag{7.7}\\
& \leq C s^{\lambda} .
\end{align*}
$$

Finally, for $I V$, observe Green's identity gives

$$
\int_{|\bar{x}-z|>s} \partial_{z_{m}} k_{l, j}(x, x-z) d z=\int_{|\bar{x}-z|=s} \nu_{m} k_{l, j}(x, x-z) d z
$$

and using (7.2), we have

$$
\left|\int_{|\bar{x}-z|=d} \nu_{m} k_{l, j}(x, x-z) d z\right| \leq C
$$

Thus $I V \leq C s^{\lambda}$.
Now we can recover curl $W$.
Theorem 7.3. Suppose the conditions in Theorem 1.3 holds; then we can construct $\widetilde{W} \in C^{\lambda}$ and compactly supported with curl $\widetilde{W}=\operatorname{curl} W$ and $\left.\widetilde{W}\right|_{\partial \Omega}=0$.

Proof. As outlined earlier, $\widetilde{W}=J_{2}(d W)$, the Theorem follows from the mapping properties of $J_{2}$ proved above.

Next we show how to recover $p$. Let $\xi \in \mathbf{R}^{n} \backslash\{0\}$ and $\gamma_{1}, \gamma_{2}$ be two unit vectors with $\left\{\xi, \gamma_{1}, \gamma_{2}\right\}$ an orthogonal set. For $s=1 / h$ we define

$$
\begin{align*}
\zeta_{1} & =-\frac{\xi}{2}+s \sqrt{1-\frac{|\xi|^{2}}{4 s^{2}}} \gamma_{1}+i s \gamma_{2}  \tag{7.8}\\
\zeta_{2} & =-\frac{\xi}{2}-s \sqrt{1-\frac{|\xi|^{2}}{4 s^{2}}} \gamma_{1}-i s \gamma_{2}
\end{align*}
$$

Let $u_{\zeta}=e^{i \zeta \cdot x}\left(\omega_{0}+\omega\right)$ be the solution of $H_{W, p} u_{\zeta}=0$ with

$$
\begin{align*}
& \omega_{0}=e^{i \chi_{\zeta} \phi^{\sharp}} \\
& \chi_{\zeta}(x)=\chi\left(x /|\zeta|^{\theta}\right),  \tag{7.9}\\
& \phi^{\sharp}(x)=N_{\mu}^{-1}\left(-\mu \cdot W^{\sharp}\right),
\end{align*}
$$

where the function $\chi \in C_{0}^{\infty}\left(\mathbf{R}^{n}\right)$, $\chi=1$ in $B(0, M / 2)$, $\chi=0$ outside $B(0, M)$, and $\bar{\Omega} \in B(0, M / 2),\left|\partial_{x}^{\alpha} W^{\sharp}(x)\right|<C h^{-|\alpha| \sigma}$. Note that this $\sigma$ is not the $\sigma$ we mentioned before. In Chapter 3, $\sigma=\sigma_{0}+\theta$, but here, we choose $\sigma$ independent of $\sigma_{0}$. Let $\mu_{j}=\zeta_{j} / s=\mu_{j}^{1}+\mu_{j}^{2}, j=1,2$. We show the following lemma.

Lemma 7.4. Suppose $W \in C^{\lambda}$ and $p=p_{1}+\operatorname{div} p_{2}, p_{1}, p_{2} \in C^{\tilde{\lambda}}$, with $\lambda>1 / 2$ and $\tilde{\lambda}>0$. Then there exists a sequence pair of $\zeta_{j}=h^{-1} \mu_{j}, j=1,2$, defined in 7.8), such that $\left\|\omega_{1}\right\|_{\dot{X}_{\mu_{j}, h}^{1 / 2}}=o\left(h^{1-\sigma-\varepsilon}\right)$, and $\left\|\omega_{2}\right\|_{\dot{X}_{\mu_{j}, h}^{1 / 2}}=o\left(h^{1-\sigma-\varepsilon}\right)$ as $h \rightarrow 0$, where $\sigma=\frac{1}{2(1+\lambda)}$, and $\varepsilon>0$ can be arbitrarily small and independent of $h, \sigma$.

Proof. Substituting $u_{\zeta}$ in $H_{W, p} u_{\zeta}=0$, we get

$$
\begin{equation*}
\left(\Delta_{\zeta}+2 W \cdot D_{\zeta}+G\right) \omega=-f \tag{7.10}
\end{equation*}
$$

or equivalently

$$
\begin{equation*}
\left(Q+2 h W \cdot(h D+\mu)+h^{2} G\right) \omega=-h^{2} f \tag{7.11}
\end{equation*}
$$

with $G=W \cdot W+D \cdot W+p$, and

$$
\begin{align*}
f= & \left(\Delta_{\zeta}+2 W \cdot D_{\zeta}+G\right) \omega_{0} \\
= & e^{i \chi_{\zeta} \phi^{\sharp}}\left[i \chi_{\zeta} \Delta \phi^{\sharp}+2 i D \chi_{\zeta} \cdot D \phi^{\sharp}+i\left(\Delta \chi_{\zeta}\right) \phi^{\sharp}+\left(\chi_{\zeta} \nabla \phi^{\sharp}+\left(\nabla \chi_{\zeta}\right) \phi^{\sharp}\right)^{2}\right. \\
& +2 \zeta \cdot\left(\nabla \chi_{\zeta}\right) \phi^{\sharp}+2 \zeta \cdot\left(\chi_{\zeta} \nabla \phi^{\sharp}\right)+2 W \cdot\left(\nabla \chi_{\zeta}\right) \phi^{\sharp}+2 W \cdot\left(\chi_{\zeta} \nabla \phi^{\sharp}\right)  \tag{7.12}\\
& \left.+2 W^{\sharp} \cdot \zeta+2 W^{b} \cdot \zeta+G\right] .
\end{align*}
$$

Since

$$
\begin{equation*}
2 \zeta \cdot \nabla \phi^{\sharp}+2 W^{\sharp} \cdot \zeta=0 \tag{7.13}
\end{equation*}
$$

which removes two terms from $f$, we have

$$
\begin{align*}
f= & \left(\Delta_{\zeta}+2 W \cdot D_{\zeta}+G\right) \omega_{0} \\
= & e^{i \chi_{\zeta} \phi^{\sharp}}\left[i \chi_{\zeta} \Delta \phi^{\sharp}+2 i D \chi_{\zeta} \cdot D \phi^{\sharp}+i\left(\Delta \chi_{\zeta}\right) \phi^{\sharp}\right. \\
& +\left(\chi_{\zeta} \nabla \phi^{\sharp}+\left(\nabla \chi_{\zeta}\right) \phi^{\sharp}\right)^{2}+2 \zeta \cdot\left(\nabla \chi_{\zeta}\right) \phi^{\sharp}  \tag{7.14}\\
& \left.+2 W \cdot\left(\nabla \chi_{\zeta}\right) \phi^{\sharp}+2 W \cdot\left(\chi_{\zeta} \nabla \phi^{\sharp}\right)+2 W^{b} \cdot \zeta+G\right] .
\end{align*}
$$

Then

$$
\begin{align*}
&\|\omega\|_{\dot{X}_{\mu_{j}, h}^{1 / 2}} \leq C h^{2-\sigma_{0}}\left\|\langle x\rangle^{1 / 2+\delta} f\right\|_{X_{\mu_{j}, h}^{-1 / 2}} \\
& \leq C h^{2-\sigma_{0}} \|\langle x\rangle^{1 / 2+\delta} e^{i \chi_{\zeta} \phi^{\sharp}}\left[i \chi_{\zeta} \Delta \phi^{\sharp}+2 i D \chi_{\zeta} \cdot D \phi^{\sharp}+i\left(\Delta \chi_{\zeta}\right) \phi^{\sharp}\right.  \tag{7.15}\\
& \quad+\left(\chi_{\zeta} \nabla \phi^{\sharp}+\left(\nabla \chi_{\zeta}\right) \phi^{\sharp}\right)^{2}+2 \zeta \cdot\left(\nabla \chi_{\zeta}\right) \phi^{\sharp} \\
&\left.\quad+2 W \cdot\left(\nabla \chi_{\zeta}\right) \phi^{\sharp}+2 W \cdot\left(\chi_{\zeta} \nabla \phi^{\sharp}\right)+2 W^{\text {b }} \cdot \zeta+G\right] \|_{X_{\mu_{j}, h}^{-1 / 2}} .
\end{align*}
$$

Next we estimate each term, since $\left|\partial^{\alpha} \phi^{\sharp}(x)\right| \leq C h^{-\sigma|\alpha|}\left\langle x_{T}\right\rangle^{-1} \chi_{B(0, M)}\left(x_{\perp}\right)$, where $x_{T}$ is the projection of $x$ to $\operatorname{span}\left\{\gamma_{1}, \gamma_{2}\right\}$ and $x_{\perp}=x-x_{T}$. Then for the first term

$$
\begin{align*}
& h^{2-\sigma_{0}}\left\|\langle x\rangle^{1 / 2+\delta} e^{i \chi_{\zeta} \phi^{\sharp}} i \chi_{\zeta} \Delta \phi^{\sharp}\right\|_{X_{\mu_{j}, h}^{-1 / 2}} \\
& \leq C h^{3 / 2-\sigma_{0}}\left\|\langle x\rangle^{1 / 2+\delta} e^{i \chi_{\zeta} \phi^{\sharp}} i \chi_{\zeta} \Delta \phi^{\sharp}\right\|_{L^{2}} \\
& \leq C h^{3 / 2-\sigma_{0}}\left(\int_{\mathbf{R}^{n}}\langle x\rangle^{2(1 / 2+\delta)}\left(\chi_{\zeta}(x)\right)^{2}\left|\Delta \phi^{\sharp}(x)\right|^{2} d x\right)^{1 / 2} \\
& \leq C h^{3 / 2-\sigma_{0}} h^{-2 \sigma}\left(\int_{\left|x_{T}\right| \leq M h^{-\theta},\left|x_{\perp}\right| \leq M}\langle x\rangle^{2(1 / 2+\delta)}\left\langle x_{T}\right\rangle^{-2} d x\right)^{1 / 2} \\
& \leq C h^{3 / 2-\sigma_{0}-2 \sigma}\left(\int_{\left|x_{T}\right| \leq M h^{-\theta}}\left\langle x_{T}\right\rangle^{-1+2 \delta} d x_{T}\right)^{1 / 2} \\
& \leq C h^{3 / 2-\sigma_{0}-2 \sigma-\theta(1 / 2+\delta)} . \tag{7.16}
\end{align*}
$$

This term has the worst behavior among the first four terms of (7.15) since the derivatives of $\chi_{\zeta}$ bring decay and the other terms only have first derivatives of $\phi^{\sharp}$. For the fifth term, we have

$$
\begin{align*}
& h^{2-\sigma_{0}}\left\|\langle x\rangle^{1 / 2+\delta} e^{i \chi_{\zeta} \phi^{\sharp}} 2 \zeta \cdot\left(\nabla \chi_{\zeta}\right) \phi^{\sharp}\right\|_{X_{\mu_{j}, h}^{-1 / 2}} \\
& \leq C h^{2-\sigma_{0}} h^{-1+\theta} h^{-1 / 2}\left\|\langle x\rangle^{1 / 2+\delta} e^{i \chi_{\zeta} \phi^{\sharp}} \phi^{\sharp}\right\|_{L^{2}} \\
& \leq C h^{1 / 2-\sigma_{0}+\theta}\left\|\langle x\rangle^{1 / 2+\delta} \phi^{\sharp}\right\|_{L^{2}}  \tag{7.17}\\
& \leq C h^{1 / 2-\sigma_{0}+\theta} h^{-\theta(1 / 2+\delta)} \\
& \leq C h^{1 / 2-\sigma_{0}+\theta(1 / 2-\delta)} .
\end{align*}
$$

The estimates for the sixth, seventh, and eighth terms are $C h^{3 / 2-\sigma_{0}+\theta}, C h^{3 / 2-\sigma_{0}-\sigma}$, and $C h^{1 / 2-\sigma_{0}+\sigma \lambda}$, respectively. For the ninth term, by Haberman and Tataru's paper [10, Lemma 3.1, p.10], we can find a sequence $\left(\mu_{j}, h_{j}\right)$ with $h_{j}$ tending to zero so that

$$
\|G\|_{X_{\mu_{j}, h_{j}}^{-1 / 2}}=o\left(h^{-1+\frac{1}{2} \tilde{\lambda}}\right) .
$$

So the estimate of the ninth term is $C h^{1-\sigma_{0}-\frac{\sigma}{2}+\frac{1}{2} \tilde{\lambda}}$.
Now consider the worst three terms: $h^{3 / 2-\sigma_{0}-2 \sigma-\theta(1 / 2+\delta)}, h^{1 / 2-\sigma_{0}+\theta(1 / 2-\delta)}$ and $h^{1 / 2-\sigma_{0}+\sigma \lambda}$. Let $3 / 2-\sigma_{0}-2 \sigma-\theta(1 / 2+\delta)=1 / 2-\sigma_{0}+\theta(1 / 2-\delta)$. We find that $\theta=1-2 \sigma$ and $h^{3 / 2-\sigma_{0}-2 \sigma-\theta(1 / 2+\delta)}=h^{1 / 2-\sigma_{0}+\theta(1 / 2-\delta)}=h^{1-\sigma-\sigma_{0}-\delta(1-2 \sigma)}$, since $\sigma_{0}, \delta$ can be arbitrarily small. Now we let $1-\sigma=1 / 2-+\sigma \lambda$; we have $\sigma=\frac{1}{2(1+\lambda)}$. Thus if we make a summary, choosing $\theta=1-2 \sigma$ and $\sigma=\frac{1}{2(1+\lambda)}$, then for any $\varepsilon>0$,

$$
\|\omega\|_{\dot{X}_{\mu_{j}, h}^{1 / 2}} \leq C h^{1-\sigma-\varepsilon} .
$$

Let $u_{\zeta_{1}}, u_{\zeta_{2}}$ satisfy $H_{W, p} u_{\zeta_{1}}=0, H_{-W, 0} u_{\zeta_{2}}=0$, respectively, and take the forms

$$
\begin{aligned}
& u_{\zeta_{1}}=e^{i \zeta_{1} \cdot x}\left(e^{i \phi_{1}^{\sharp}}+\omega_{1}\right) \\
& u_{\zeta_{2}}=e^{i \zeta_{2} \cdot x}\left(e^{-i \phi_{2}^{\sharp}}+\omega_{2}\right)
\end{aligned}
$$

Define a scattering transform

$$
\tilde{t}(\xi)=\left\langle\left(\Lambda_{W, p}-\Lambda_{-W, 0}\right)\left(u_{\zeta_{1}} \mid \partial \Omega\right), v_{\zeta_{2}} \mid \partial \Omega\right\rangle=\left\langle p u_{\zeta_{1}}, u_{\zeta_{2}}\right\rangle
$$

Substituting $u_{\zeta_{1}}$ and $u_{\zeta_{2}}$, we get

$$
\begin{aligned}
\tilde{t}(\xi) & =\left\langle e^{-i x \cdot \xi} p\left(e^{i \phi_{1}^{\sharp}}+\omega_{1}\right),\left(e^{-i \phi_{2}^{\#}}+\omega_{2}\right)\right\rangle \\
& =\left\langle e^{-i x \cdot \xi} p e^{i \phi_{1}^{\sharp}}, e^{-i \phi_{2}^{\#}}\right\rangle+\left\langle e^{-i x \cdot \xi} p e^{i \phi_{1}^{\sharp}}, \omega_{2}\right\rangle+\left\langle e^{-i x \cdot \xi} p \omega_{1}, e^{-i \phi_{2}^{\sharp}}\right\rangle+\left\langle e^{-i x \cdot \xi} p \omega_{1}, \omega_{2}\right\rangle
\end{aligned}
$$

For the first term, since

$$
\zeta_{1} \cdot \nabla \phi_{1}^{\sharp}+W^{\sharp} \cdot \zeta_{1}=0,
$$

$$
\zeta_{2} \cdot \nabla \phi_{2}^{\sharp}-W^{\sharp} \cdot \zeta_{2}=0
$$

and $\zeta_{1}+\zeta_{2}=-\xi=s\left(\mu_{1}+\mu_{2}\right)$ as $s \rightarrow \infty$ or $\mu_{1}=-\mu_{2}-h \xi$. Furthermore,

$$
\phi_{1}^{\sharp}=N_{\mu_{1}}^{-1}\left(\mu_{1} \cdot W^{\sharp}\right)=\frac{1}{2 \pi} \int_{\mathbf{R}^{2}} \frac{1}{y_{1}+i y_{2}}\left(\mu_{1} \cdot W^{\sharp}\right)\left(x-y_{1} \mu_{1}^{1}-y_{2} \mu_{1}^{2}\right) d y_{1} d y_{2}
$$

and

$$
\phi_{2}^{\sharp}=N_{\mu_{2}}^{-1}\left(-\mu_{2} \cdot W^{\sharp}\right)=\frac{1}{2 \pi} \int_{\mathbf{R}^{2}} \frac{1}{y_{1}+i y_{2}}\left(-\mu_{2} \cdot W^{\sharp}\right)\left(x+y_{1} \mu_{2}^{1}+y_{2} \mu_{2}^{2}\right) d y_{1} d y_{2} .
$$

So

$$
\begin{align*}
& \phi_{1}^{\sharp}-\phi_{2}^{\sharp}=\frac{1}{2 \pi} \int_{\mathbf{R}^{2}} \frac{1}{y_{1}+i y_{2}}\left(\mu_{1} \cdot W^{\sharp}\right)\left(x-y_{1} \mu_{1}^{1}-y_{2} \mu_{1}^{2}\right) d y_{1} d y_{2} \\
&+\frac{1}{2 \pi} \int_{\mathbf{R}^{2}} \frac{1}{y_{1}+i y_{2}}\left(\mu_{2} \cdot W^{\sharp}\right)\left(x+y_{1} \mu_{2}^{1}+y_{2} \mu_{2}^{2}\right) d y_{1} d y_{2} \\
&=\frac{1}{2 \pi} \int_{\mathbf{R}^{2}} \frac{1}{y_{1}+i y_{2}}\left(\left(-\mu_{2}-h \xi\right) \cdot W^{\sharp}\right)\left(x-y_{1} \mu_{1}^{1}-y_{2} \mu_{1}^{2}\right) d y_{1} d y_{2} \\
&+\frac{1}{2 \pi} \int_{\mathbf{R}^{2}} \frac{1}{y_{1}+i y_{2}}\left(\mu_{2} \cdot W^{\sharp}\right)\left(x+y_{1} \mu_{2}^{1}+2 \mu_{2}^{2}\right) d y_{1} d y_{2} \\
&=\frac{1}{2 \pi} \int_{\mathbf{R}^{2}} \frac{1}{y_{1}+i y_{2}}\left((-h \xi) \cdot W^{\sharp}\right)\left(x-y_{1} \mu_{1}^{1}-y_{2} \mu_{1}^{2}\right) d y_{1} d y_{2} \\
&+\frac{1}{2 \pi} \int_{\mathbf{R}^{2}} \frac{\left(\mu_{2} \cdot W^{\sharp}\right)\left(x+y_{1} \mu_{2}^{1}+y_{2} \mu_{2}^{2}\right)-\left(\mu_{2} \cdot W^{\sharp}\right)\left(x-y_{1} \mu_{1}^{1}-y_{2} \mu_{1}^{2}\right)}{y_{1}+i y_{2}} d y_{1} d y_{2} \\
&=I+I I . \tag{7.18}
\end{align*}
$$

Because $\partial_{x}^{\alpha} W^{\sharp} \leq C h^{-\sigma|\alpha|}$ and $\left|\left(x+y_{1} \mu_{2}^{1}+y_{2} \mu_{2}^{2}\right)-\left(x-y_{1} \mu_{1}^{1}-y_{2} \mu_{1}^{2}\right)\right| \leq C h\left|y_{1}\right|$, then $I \leq C h$ and $I I \leq C h^{1-\sigma}$, similarly, we can prove $\left|\nabla\left(\phi_{1}^{\sharp}-\phi_{2}^{\sharp}\right)\right| \leq C h^{1-2 \sigma}$.

It follows that

$$
\left\langle e^{-i x \cdot \xi} p e^{i \phi_{1}^{\#}}, e^{-i \phi_{2}^{\#}}\right\rangle \rightarrow \hat{p}(\xi)
$$

as $s \rightarrow \infty$.
Next the second term: if we consider $e^{i \phi_{1}^{\sharp}}$ as a symbol, from Lemma 3.2, we have

$$
\begin{align*}
\mid\left\langle e^{-i x \cdot \xi} p e^{i \phi_{1}^{\sharp}}, \omega_{2}\right\rangle & \leq C\left\|e^{-i x \cdot \xi} p e^{i \phi_{1}^{\sharp}}\right\|_{\dot{X}_{\mu, h}^{-1 / 2}}\left\|\omega_{2}\right\|_{\dot{X}_{\mu, h}^{1 / 2}} \\
& \leq C h^{-\sigma / 2}\left\|e^{-i x \cdot \xi} p\right\|_{\dot{X}_{\mu, h}^{-1 / 2}}\left\|\omega_{2}\right\|_{\dot{X}_{\mu, h}^{1 / 2}} \\
& \leq C h^{-\sigma / 2} h^{-1+\frac{1}{2} \tilde{\lambda}} h^{1-\sigma-\varepsilon}  \tag{7.19}\\
& \leq C h^{\frac{1}{2} \tilde{\lambda}-\frac{3}{2} \sigma-\varepsilon} \\
& \leq C h^{1 / 2\left(\tilde{\lambda}-\frac{3}{2(1+\lambda)}\right)-\varepsilon}
\end{align*}
$$

where the estimate of $\left\|e^{-i x \cdot \xi} p e^{i \phi_{1}^{\sharp}}\right\|_{\dot{X}_{\mu, h}^{-1 / 2}}$ leads to an average estimate found in Haberman and Tataru [10] [Lemma 3.1]. The third term is similar to the second.

Finally we look at the last term,

$$
\begin{align*}
\left|\left\langle e^{-i x \cdot \xi} p \omega_{1}, \omega_{2}\right\rangle\right| & \leq C h^{-2+\tilde{\lambda}}\left\|\omega_{1}\right\|_{\dot{X}_{\mu, h}^{1 / 2}}\left\|\omega_{2}\right\|_{\dot{X}_{\mu, h}^{1 / 2}} \\
& \leq C h^{-2+\tilde{\lambda}} h^{1-\sigma-\varepsilon} h^{1-\sigma-\varepsilon}  \tag{7.20}\\
& \leq C h^{\tilde{\lambda}-\frac{1}{1+\lambda}-2 \varepsilon}
\end{align*}
$$

where the estimate of the operator $e^{-i x \cdot \xi} p$ is follow to Haberman and Tataru [10][Theorem 2.1].

Above all, if $\tilde{\lambda}-\frac{3}{2(1+\lambda)}>0$, or equivalently $\tilde{\lambda}(1+\lambda)>\frac{3}{2}$, we can recover $\hat{p}$, which is the Fourier transform of $p$. Then we could construct our electric potential $p$. If we let $\tilde{\lambda}$ approaches 1 , then $\lambda$ can approaches to $1 / 2$ from above, and if we let $\lambda$ approaches to $1, \tilde{\lambda}$ converges to $3 / 4$ from above. To make a summary, we can say that for any $\lambda \in(1 / 2,1), \tilde{\lambda} \in(3 / 4,1)$ and $\tilde{\lambda}(1+\lambda)>\frac{3}{2}, p$ can be determined by DN map.

## Chapter 8 Future work

In the previous chapter, we recovered $p$ by using CGO solutions of the form $u=$ $e^{i x \cdot \zeta}\left(e^{i \phi^{\sharp}}+\omega\right)$. The estimate 7.19 requires $\tilde{\lambda}-\frac{3}{2(1+\lambda)}>0$, the estimate 7.20 requires $\tilde{\lambda}-\frac{1}{1+\lambda}>0$. We do not expect this result to be sharp and in future work we would like to recover $p$ for a larger range of $\lambda$ and $\tilde{\lambda}$.

Other interesting questions include studying the stability of the recovery process for non-smooth potentials. This would extend the work of Leo Tzou [25]. We can also consider extending Haberman's methods for magnetic Schrödinger operators to the case when the magnetic potential $W$ is not small. Finally, there is much interesting work for the inverse boundary value problem in the case when we only have data on part of the boundary. The inverse boundary value problem for the magnetic Schrödinger operator was studied by Chung [6]. Another area of investigation is to study this partial data problem when the potentials are not smooth.

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Vita

## Vita

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