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# On stability estimates in the Gel'fand-Calderon inverse problem 

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Abstract. We prove new global stability estimates for the Gel'fand-Calderon inverse problem in 3D.

## 1. Introduction

We consider the equation

$$
\begin{equation*}
-\Delta \psi+v(x) \psi=0, \quad x \in D \tag{1.1}
\end{equation*}
$$

where

$$
\begin{equation*}
D \text { is an open bounded domain in } \mathbb{R}^{d}, d \geq 2, \partial D \in C^{2}, v \in L^{\infty}(D) . \tag{1.2}
\end{equation*}
$$

Equation (1.1) arises, in particular, in quantum mechanics, acoustics, electrodynamics. Formally, (1.1) looks as the Schrödinger equation with potential $v$ at zero energy.

We consider the map $\Phi$ such that

$$
\begin{equation*}
\left.\frac{\partial \psi}{\partial \nu}\right|_{\partial D}=\Phi\left(\left.\psi\right|_{\partial D}\right) \tag{1.3}
\end{equation*}
$$

for all sufficiently regular solutions $\psi$ of (1.1) in $\bar{D}=D \cup \partial D$, where $\nu$ is the outward normal to $\partial D$. Here we assume also that

$$
\begin{equation*}
0 \text { is not a Dirichlet eigenvalue for the operator }-\Delta+v \text { in } D \text {. } \tag{1.4}
\end{equation*}
$$

The map $\Phi$ is called the Dirichlet-to-Neumann map for equation (1.1) and is considered as boundary measurements for (physical model described by) (1.1).

We consider the following inverse boundary value problem for equation (1.1):
Problem 1.1. Given $\Phi$, find $v$.
This problem can be considered as the Gel'fand inverse boundary value problem for the Schrödinger equation at zero energy (see [9], [16]). This problem can be also considered as a generalization of the Calderon problem of the electrical impedance tomography (see [5], [23], [16]).

Concerning results given in the literature on Problem 1.1 (in its Calderon or Gel'fand form) see [6], [11], [23], [10] (note added in proof), [16], [1], [14], [15], [3], [22], [13], [17], [19], [12], [4], [21], [20] and references therein.

In particular, in [21] it was shown that the Alessandrini stability estimates of [1] for Problem 1.1 in dimension $d \geq 3$ (see Theorem 2.1 of the next section) admit some principle improvement. These new stability estimates (see Theorem 2.2 of the next section) were found in [21] using methods developed in [17], [18], [19]. These methods include, in particular: (1) the $\bar{\partial}$-approach to inverse "scattering" at zero energy in dimension $d \geq 3$,
going back to [2], [10], and (2) the reduction of Problem 1.1 to inverse "scattering" at zero energy, going back to [16].

However, a complete proof of the aforementioned new stability estimates for Problem 1.1 in dimension $d \geq 2$ was given in [21] in the Born approximation (that is in the linear approximation near zero potential) only. Besides, a scheme of proof of these estimates was also mentioned in [21] for potentials with sufficiently small norm in dimension $d=3$. (In this scheme [21] refers, in particular, to results of [19].)

In the present work we give a complete proof of these new stability estimates (Theorem 2.2 of the next section) in the general (or by other words global) case in dimension $d=3$. In this proof we use, in particular, results of the recent work [20].

## 2. Stability estimates

As in [21] we assume for simplicity that

$$
\begin{align*}
& D \text { is an open bounded domain in } \mathbb{R}^{d}, \partial D \in C^{2},  \tag{2.1}\\
& v \in W^{m, 1}\left(\mathbb{R}^{d}\right) \text { for some } m>d, \text { supp } v \subset D, d \geq 2,
\end{align*}
$$

where

$$
\begin{equation*}
W^{m, 1}\left(\mathbb{R}^{d}\right)=\left\{v: \partial^{J} v \in L^{1}\left(\mathbb{R}^{d}\right),|J| \leq m\right\}, \quad m \in \mathbb{N} \cup 0, \tag{2.2}
\end{equation*}
$$

where

$$
J \in(\mathbb{N} \cup 0)^{d}, \quad|J|=\sum_{i=1}^{d} J_{i}, \quad \partial^{J} v(x)=\frac{\partial^{|J|} v(x)}{\partial x_{1}^{J_{1}} \ldots \partial x_{d}^{J_{d}}}
$$

Let

$$
\begin{equation*}
\|v\|_{m, 1}=\max _{|J| \leq m}\left\|\partial^{J} v\right\|_{L^{1}\left(\mathbb{R}^{d}\right)} \tag{2.3}
\end{equation*}
$$

Let

$$
\begin{align*}
& \|A\| \text { denote the norm of an operator } \\
& A: L^{\infty}(\partial D) \rightarrow L^{\infty}(\partial D) \tag{2.4}
\end{align*}
$$

We recall that if $v_{1}, v_{2}$ are potentials satisfying (1.2), (1.3), where $D$ is fixed, then

$$
\begin{equation*}
\Phi_{1}-\Phi_{2} \text { is a compact operator in } L^{\infty}(\partial D) \tag{2.5}
\end{equation*}
$$

where $\Phi_{1}, \Phi_{2}$ are the DtN maps for $v_{1}, v_{2}$ respectively, see [16], [17]. Note also that (2.1) $\Rightarrow$ (1.2).

Theorem 2.1 (variation of the result of [1]). Let conditions (1.4), (2.1) hold for potentials $v_{1}$ and $v_{2}$, where $D$ is fixed, $d \geq 3$. Let $\left\|v_{j}\right\|_{m, 1} \leq N, j=1,2$, for some $N>0$. Let $\Phi_{1}, \Phi_{2}$ denote the DtN maps for $v_{1}, v_{2}$, respectively. Then

$$
\begin{equation*}
\left\|v_{1}-v_{2}\right\|_{L^{\infty}(D)} \leq c_{1}\left(\ln \left(1+\left\|\Phi_{1}-\Phi_{2}\right\|^{-1}\right)\right)^{-\alpha_{1}} \tag{2.6}
\end{equation*}
$$

where $c_{1}=c_{1}(N, D, m), \alpha_{1}=(m-d) / m,\left\|\Phi_{1}-\Phi_{2}\right\|$ is defined according to (2.4).
As it was mentioned in [21], Theorem 2.1 follows from formulas (3.9)-(3.11), (4.1) (of Sections 3 and 4).

A disadvantage of estimate (2.6) is that

$$
\begin{equation*}
\alpha_{1}<1 \text { for any } m>d \text { even if } m \text { is very great. } \tag{2.7}
\end{equation*}
$$

Theorem 2.2. Let the assumptions of Theorem 2.1 hold. Then

$$
\begin{equation*}
\left\|v_{1}-v_{2}\right\|_{L^{\infty}(D)} \leq c_{2}\left(\ln \left(1+\left\|\Phi_{1}-\Phi_{2}\right\|^{-1}\right)\right)^{-\alpha_{2}} \tag{2.8}
\end{equation*}
$$

where $c_{2}=c_{2}(N, D, m), \alpha_{2}=m-d,\left\|\Phi_{1}-\Phi_{2}\right\|$ is defined according to (2.4).
A principal advantage of estimate (2.8) in comparison with (2.6) is that

$$
\begin{equation*}
\alpha_{2} \rightarrow+\infty \text { as } m \rightarrow+\infty, \tag{2.9}
\end{equation*}
$$

in contrast with (2.7).
In the Born approximation, that is in the linear approximation near zero potential, Theorem 2.2 was proved in [21].

For sufficiently small $N$ in dimension $d=3$, a scheme of proof of Theorem 2.2 was also mentioned in [21]. This scheme involves, in particular, results of [17], [19].

In the general (or by other words global) case Theorem 2.2 in dimension $d=3$ is proved in Section 7. This proof involves, in particular, results of [17], [20].

## 3. Faddeev functions

We consider the Faddeev functions $G, \psi$ and $h$ (see [7], [8], [10], [16]):

$$
\begin{align*}
& \psi(x, k)=e^{i k x}+\int_{\mathbb{R}^{d}} G(x-y, k) v(y) \psi(y, k) d y  \tag{3.1}\\
& G(x, k)=e^{i k x} g(x, k), \quad g(x, k)=-(2 \pi)^{-d} \int_{\mathbb{R}^{d}} \frac{e^{i \xi x} d \xi}{\xi^{2}+2 k \xi} \tag{3.2}
\end{align*}
$$

where $x \in \mathbb{R}^{d}, k \in \Sigma$,

$$
\begin{align*}
& \Sigma=\left\{k \in \mathbb{C}^{d}: \quad k^{2}=k_{1}^{2}+\ldots+k_{d}^{2}=0\right\}  \tag{3.3}\\
& h(k, l)=(2 \pi)^{-d} \int_{\mathbb{R}^{d}} e^{-i l x} v(x) \psi(x, k) d x \tag{3.4}
\end{align*}
$$

where $(k, l) \in \Theta$,

$$
\begin{equation*}
\Theta=\{k \in \Sigma, \quad l \in \Sigma: \quad \operatorname{Im} k=\operatorname{Im} l\} . \tag{3.5}
\end{equation*}
$$

One can consider (3.1), (3.4) assuming that $v$ is a sufficiently regular function on $\mathbb{R}^{d}$ with sufficient decay at infinity. For example, one can consider (3.1), (3.4) assuming that (1.2) holds.

We recall that:

$$
\begin{equation*}
\Delta G(x, k)=\delta(x), \quad x \in \mathbb{R}^{d}, \quad k \in \Sigma \tag{3.6}
\end{equation*}
$$

formula (3.1) at fixed $k$ is considered as an equation for

$$
\begin{equation*}
\psi=e^{i k x} \mu(x, k) \tag{3.7}
\end{equation*}
$$

where $\mu$ is sought in $L^{\infty}\left(\mathbb{R}^{d}\right)$; as a corollary of (3.1),(3.2), (3.6), $\psi$ satisfies (1.1); $h$ of (3.4) is a generalized "scattering" amplitude in the complex domain at zero energy.

Note that, actually, $G, \psi, h$ of (3.1)-(3.5) are zero energy restrictions of functions introduced by Faddeev as extentions to the complex domain of some functions of the classical scattering theory for the Schrödinger equation at positive energies. In addition, $G, \psi, h$ in their zero energy restriction were considered for the first time in [2]. The Faddeev functions $G, \psi, h$ were, actually, rediscovered in [2].

We recall also that, under the assumptions of Theorem 2.1,

$$
\begin{equation*}
\mu(x, k) \rightarrow 1 \text { as }|I m k| \rightarrow \infty \quad \text { (uniformly in } x \text { ) } \tag{3.8}
\end{equation*}
$$

and, for any $\sigma>1$,

$$
\begin{equation*}
|\mu(x, k)|<\sigma \text { for }|I m k| \geq r_{1}(N, D, m, \sigma) \tag{3.9}
\end{equation*}
$$

where $x \in \mathbb{R}^{d}, k \in \Sigma$;

$$
\begin{gather*}
\hat{v}(p)=\lim _{\substack{(k, l) \in \Theta, k-l=p \\
|I m k|=|\operatorname{Im} m| \rightarrow \infty}} h(k, l) \text { for any } p \in \mathbb{R}^{d},  \tag{3.10}\\
|\hat{v}(p)-h(k, l)| \leq \frac{c_{3}(D, m) N^{2}}{\rho} \text { for }(k, l) \in \Theta, p=k-l,  \tag{3.11}\\
|\operatorname{Im} k|=|\operatorname{Im} l|=\rho \geq r_{2}(N, D, m),
\end{gather*}
$$

where

$$
\begin{equation*}
\hat{v}(p)=\left(\frac{1}{2 \pi}\right)^{d} \int_{\mathbb{R}^{d}} e^{i p x} v(x) d x, \quad p \in \mathbb{R}^{d} \tag{3.12}
\end{equation*}
$$

Results of the type (3.8), (3.9) go back to [2]. Results of the type (3.10), (3.11) (with less precise right-hand side in (3.11)) go back to [10]. Estimates (3.8), (3.11) are related also with some important $L_{2}$-estimate going back to [23] on the Green function $g$ of (3.1).

Note also that in some considerations it is convenient to consider $h$ on $\Theta$ as $H$ on $\Omega$, where

$$
\begin{align*}
& h(k, l)=H(k, k-l), \quad(k, l) \in \Theta,  \tag{3.13}\\
& H(k, p)=h(k, k-p), \quad(k, p) \in \Omega, \\
\Omega= & \left\{k \in \mathbb{C}^{d}, p \in \mathbb{R}^{d}: k^{2}=0, p^{2}=2 k p\right\} . \tag{3.14}
\end{align*}
$$

For more information on properties of the Faddeev functions $G, \psi, h$, see [10], [17], [20] and references therein.

In the next section we recall that Problem 1.1 (of Introduction) admits a reduction to the following inverse "scattering" problem:

Problem 3.1. Given $h$ on $\Theta$, find $v$ on $\mathbb{R}^{d}$.
4. Reduction of [16], [17]

Let conditions (1.2), (1.4) hold for potentials $v_{1}$ and $v_{2}$, where $D$ is fixed. Let $\Phi_{i}, \psi_{i}$, $h_{i}$ denote the DtN map $\Phi$ and the Faddeev functions $\psi, h$ for $v=v_{i}, i=1,2$. Let also
$\Phi_{i}(x, y)$ denote the Schwartz kernel $\Phi(x, y)$ of the integral operator $\Phi$ for $v=v_{i}, i=1,2$. Then (see [17] for details):

$$
\begin{equation*}
h_{2}(k, l)-h_{1}(k, l)=\left(\frac{1}{2 \pi}\right)^{d} \int_{\partial D \partial D} \int_{1}(x,-l)\left(\Phi_{2}-\Phi_{1}\right)(x, y) \psi_{2}(y, k) d y d x \tag{4.1}
\end{equation*}
$$

where $(k, l) \in \Theta$;

$$
\begin{align*}
& \psi_{2}(x, k)=\psi_{1}(x, k)+\int_{\partial D} A(x, y, k) \psi_{2}(y, k) d y, x \in \partial D  \tag{4.2a}\\
& A(x, y, k)=\int_{\partial D} R_{1}(x, z, k)\left(\Phi_{2}-\Phi_{1}\right)(z, y) d z, x, y \in \partial D  \tag{4.2b}\\
& R_{1}(x, y, k)=G(x-y, k)+\int_{\mathbb{R}^{d}} G(x-z, k) v_{1}(z) R_{1}(z, y, k) d z, x, y \in \mathbb{R}^{d} \tag{4.3}
\end{align*}
$$

where $k \in \Sigma$. Note that: (4.1) is an explicit formula, (4.2a) is considered as an equation for finding $\psi_{2}$ on $\partial D$ from $\psi_{1}$ on $\partial D$ and $A$ on $\partial D \times \partial D$ for each fixed $k,(4.2 \mathrm{~b})$ is an explicit formula, (4.3) is an equation for finding $R_{1}$ from $G$ and $v_{1}$, where $G$ is the function of (3.2).

Note that formulas and equations (4.1)-(4.3) for $v_{1} \equiv 0$ were given in [16] (see also [10] (Note added in proof), [14], [15]). In this case $h_{1} \equiv 0, \psi_{1}=e^{i k x}, R_{1}=G(x-y, k)$. Formulas and equations (4.1)-(4.3) for the general case were given in [17].

Formulas and equations (4.1)-(4.3) with fixed background potential $v_{1}$ reduce Problem 1.1 (of Introduction) to Problem 3.1 (of Section 3).

## 5. Some considerations related with $\Theta$ and $\Omega$

5.1 Some subsets of $\Theta$ and $\Omega$. Let

$$
\begin{align*}
& \mathcal{B}_{r}=\left\{p \in \mathbb{R}^{d}:|p|<r\right\}, \quad \partial \mathcal{B}_{r}=\left\{p \in \mathbb{R}^{d}:|p|=r\right\} \\
& \overline{\mathcal{B}}_{r}=\mathcal{B}_{r} \cup \partial \mathcal{B}_{r}, \quad \text { where } r>0 \tag{5.1}
\end{align*}
$$

In addition to $\Theta$ of (3.5), we consider, in particular, the following its subsets:

$$
\begin{align*}
& \Theta_{\rho}=\{(k, l) \in \Theta:|\operatorname{Im} k|=|\operatorname{Im} l|<\rho\}, \\
& b \Theta_{\rho}=\{(k, l) \in \Theta:|\operatorname{Im} k|=|\operatorname{Im} l|=\rho\}, \\
& \bar{\Theta}_{\rho}=\Theta_{\rho} \cup b \Theta_{\rho},  \tag{5.2}\\
& \Theta_{\rho, \tau}^{\infty}=\left\{(k, l) \in \Theta \backslash \bar{\Theta}_{\rho}: k-l \in \mathcal{B}_{2 \rho \tau}\right\}, \\
& b \Theta_{\rho, \tau}=\left\{(k, l) \in b \Theta_{\rho}: k-l \in \mathcal{B}_{2 \rho \tau}\right\},
\end{align*}
$$

where $\rho>0,0<\tau<1$, and $\mathcal{B}_{r}$ is defined in (5.1).
In addition to $\Omega$ of (3.14), we consider, in particular, the following its subsets:

$$
\begin{align*}
& \Omega_{\rho}=\{(k, p) \in \Omega: \mid \text { Im } k \mid<\rho\}, \\
& b \Omega_{\rho}=\{(k, p) \in \Omega:|\operatorname{Im} k|=\rho\}, \\
& \bar{\Omega}_{\rho}=\Omega_{\rho} \cup b \Omega_{\rho},  \tag{5.3}\\
& \Omega_{\rho, \tau}^{\infty}=\left\{(k, p) \in \Omega \backslash \bar{\Omega}_{\rho}: p \in \mathcal{B}_{2 \rho \tau}\right\}, \\
& b \Omega_{\rho, \tau}=\left\{(k, p) \in b \Omega_{\rho}: p \in \mathcal{B}_{2 \rho \tau}\right\},
\end{align*}
$$

where $\rho>0,0<\tau<1$, and $\mathcal{B}_{r}$ is defined in (5.1).
Note that

$$
\begin{align*}
& \Omega \approx \Theta, \Omega_{\rho} \approx \Theta_{\rho}, b \Omega_{\rho} \approx b \Theta_{\rho} \\
& \Omega_{\rho, \tau}^{\infty} \approx \Theta_{\rho, \tau}^{\infty}, b \Omega_{\rho, \tau} \approx b \Theta_{\rho, \tau} \tag{5.4}
\end{align*}
$$

or, more precisely,

$$
\begin{align*}
& (k, p) \in \Omega \Rightarrow(k, k-p) \in \Theta, \quad(k, l) \in \Theta \Rightarrow(k, k-l) \in \Omega \\
& \text { and the same for } \Omega_{\rho}, b \Omega_{\rho}, \Omega_{\rho, \tau}^{\infty}, b \Omega_{\rho, \tau}  \tag{5.5}\\
& \text { and } \Theta_{\rho}, b \Theta_{\rho}, \Theta_{\rho, \tau}^{\infty}, b \Theta_{\rho, \tau}, \quad \text { respectively, in place of } \Omega \text { and } \Theta \text {. }
\end{align*}
$$

We consider also, in particular,

$$
\begin{align*}
& \Omega_{\nu}=\left\{(k, p) \in \Omega: p \notin \mathcal{L}_{\nu}\right\} \\
& \Omega_{\rho, \tau, \nu}^{\infty}=\Omega_{\rho, \tau}^{\infty} \cap \Omega_{\nu}, b \Omega_{\rho, \tau, \nu}=b \Omega_{\rho, \tau} \cap \Omega_{\nu} \tag{5.6}
\end{align*}
$$

where

$$
\begin{equation*}
\mathcal{L}_{\nu}=\left\{p \in \mathbb{R}^{d}: p=t \nu, t \in \mathbb{R}\right\} \tag{5.7}
\end{equation*}
$$

$\nu \in \mathbb{S}^{d-1}, \rho>0,0<\tau<1$.
5.2. Coordinates on $\Omega$ for $d=3$. In this subsection we assume that $d=3$ in formulas (3.5), (3.14), (5.1)-(5.7).

For $p \in \mathbb{R}^{3} \backslash \mathcal{L}_{\nu}$ we consider $\theta(p)$ and $\omega(p)$ such that

$$
\begin{align*}
& \theta(p), \omega(p) \text { smoothly depend on } p \in \mathbb{R}^{3} \backslash \mathcal{L}_{\nu} \\
& \text { take values in } \mathbb{S}^{2}, \text { and }  \tag{5.8}\\
& \theta(p) p=0, \omega(p) p=0, \theta(p) \omega(p)=0
\end{align*}
$$

where $\mathcal{L}_{\nu}$ is defined by (5.7) (for $d=3$ ).
Assumptions (5.8) imply that

$$
\begin{equation*}
\omega(p)=\frac{p \times \theta(p)}{|p|} \text { for } p \in \mathbb{R}^{3} \backslash \mathcal{L}_{\nu} \tag{5.9a}
\end{equation*}
$$

or

$$
\begin{equation*}
\omega(p)=-\frac{p \times \theta(p)}{|p|} \text { for } p \in \mathbb{R}^{3} \backslash \mathcal{L}_{\nu} \tag{5.9b}
\end{equation*}
$$

where $\times$ denotes vector product.
To satisfy (5.8), (5.9a) we can take

$$
\begin{equation*}
\theta(p)=\frac{\nu \times p}{|\nu \times p|}, \omega(p)=\frac{p \times \theta(p)}{|p|}, p \in \mathbb{R}^{3} \backslash \mathcal{L}_{\nu} \tag{5.10}
\end{equation*}
$$

Let $\theta, \omega$ satisfy (5.8). Then (according to [19]) the following formulas give a diffeomorphism between $\Omega_{\nu}$ and $(\mathbb{C} \backslash 0) \times\left(\mathbb{R}^{3} \backslash \mathcal{L}_{\nu}\right)$ :

$$
\begin{equation*}
(k, p) \rightarrow(\lambda, p), \quad \text { where } \lambda=\lambda(k, p)=\frac{2 k(\theta(p)+i \omega(p))}{i|p|} \tag{5.11a}
\end{equation*}
$$

$$
\begin{align*}
& (\lambda, p) \rightarrow(k, p), \quad \text { where } \quad k=k(\lambda, p)=\kappa_{1}(\lambda, p) \theta(p)+\kappa_{2}(\lambda, p) \omega(p)+\frac{p}{2}  \tag{5.11b}\\
& \kappa_{1}(\lambda, p)=\frac{i|p|}{4}\left(\lambda+\frac{1}{\lambda}\right), \quad \kappa_{2}(\lambda, p)=\frac{|p|}{4}\left(\lambda-\frac{1}{\lambda}\right)
\end{align*}
$$

where $(k, p) \in \Omega_{\nu},(\lambda, p) \in(\mathbb{C} \backslash 0) \times\left(\mathbb{R}^{3} \backslash \mathcal{L}_{\nu}\right)$. In addition, formulas (5.11a), (5.11b) for $\lambda(k)$ and $k(\lambda)$ at fixed $p \in \mathbb{R}^{3} \backslash \mathcal{L}_{\nu}$ give a diffeomorphism between $Z_{p}=\left\{k \in \mathbb{C}^{3}:(k, p) \in \Omega\right\}$ for fixed $p$ and $\mathbb{C} \backslash 0$.

In addition, for $k$ and $\lambda$ of (5.11) we have that

$$
\begin{equation*}
|\operatorname{Im} k|=\frac{|p|}{4}\left(|\lambda|+\frac{1}{|\lambda|}\right), \quad|\operatorname{Re} k|=\frac{|p|}{4}\left(|\lambda|+\frac{1}{|\lambda|}\right), \tag{5.12}
\end{equation*}
$$

where $(k, p) \in \Omega_{\nu},(\lambda, p) \in(\mathbb{C} \backslash 0) \times\left(\mathbb{R}^{3} \backslash \mathcal{L}_{\nu}\right)$.
Let

$$
\begin{gather*}
\Lambda_{\rho, \nu}=\left\{(\lambda, p): \lambda \in \mathcal{D}_{\rho /|p|}, \quad p \in \mathbb{R}^{3} \backslash \mathcal{L}_{\nu}\right\},  \tag{5.13}\\
\Lambda_{\rho, \tau, \nu}=\left\{(\lambda, p): \lambda \in \mathcal{D}_{\rho /|p|}, \quad p \in \mathbb{R}^{3} \backslash \mathcal{L}_{\nu}, \quad|p|<2 \tau \rho\right\}, \\
b \Lambda_{\rho, \tau, \nu}=\left\{(\lambda, p): \lambda \in \mathcal{T}_{\rho /|p|}, \quad p \in \mathbb{R}^{3} \backslash \mathcal{L}_{\nu}, \quad|p|<2 \tau \rho\right\},
\end{gather*}
$$

where $\rho>0,0<\tau<1, \nu \in \mathbb{S}^{2}$,

$$
\begin{align*}
& \mathcal{D}_{r}=\left\{\lambda \in \mathbb{C} \backslash 0: \frac{1}{4}\left(|\lambda|+|\lambda|^{-1}\right)>r\right\}, r>0  \tag{5.14}\\
& \mathcal{T}_{r}=\left\{\lambda \in \mathbb{C}: \frac{1}{4}\left(|\lambda|+|\lambda|^{-1}\right)=r\right\}, r \geq 1 / 2 \tag{5.15}
\end{align*}
$$

$\mathcal{L}_{\nu}$ is defined by (5.7) (for $d=3$ ).
Note that

$$
\begin{align*}
& \Lambda_{\rho, \tau, \nu}=\Lambda_{\rho, \tau, \nu}^{+} \cup \Lambda_{\rho, \tau, \nu}^{-}, \Lambda_{\rho, \tau, \nu}^{+} \cap \Lambda_{\rho, \tau, \nu}^{-}=\emptyset \\
& b \Lambda_{\rho, \tau, \nu}=b \Lambda_{\rho, \tau, \nu}^{+} \cup b \Lambda_{\rho, \tau, \nu}^{-} \tag{5.16}
\end{align*}
$$

where

$$
\begin{align*}
& \Lambda_{\rho, \tau, \nu}^{ \pm}=\left\{(\lambda, p): \quad \lambda \in \mathcal{D}_{\rho /|p|}^{ \pm}, \quad p \in \mathcal{B}_{2 \tau \rho} \backslash \mathcal{L}_{\nu}\right\},  \tag{5.17}\\
& b \Lambda_{\rho, \tau, \nu}^{ \pm}=\left\{(\lambda, p): \quad \lambda \in \mathcal{T}_{\rho /|p|}^{ \pm}, \quad p \in \mathcal{B}_{2 \tau \rho} \backslash \mathcal{L}_{\nu}\right\}, \\
& \mathcal{D}_{r}^{ \pm}=\left\{\lambda \in \mathbb{C} \backslash 0: \quad \frac{1}{4}\left(|\lambda|+|\lambda|^{-1}\right)>r, \quad|\lambda|^{ \pm 1}<1\right\}, \\
& \mathcal{T}_{r}^{ \pm}=\left\{\lambda \in \mathbb{C}: \quad \frac{1}{4}\left(|\lambda|+|\lambda|^{-1}\right)=r, \quad|\lambda|^{ \pm 1} \leq 1\right\}, r>1 / 2, \tag{5.18}
\end{align*}
$$

where $\rho>0, \tau \in] 0,1\left[, \nu \in \mathbb{S}^{2}\right.$.
Using (5.12) one can see that formulas (5.11) give also the following diffeomorphisms

$$
\begin{align*}
& \Omega_{\nu} \backslash \bar{\Omega}_{\rho} \approx \Lambda_{\rho, \nu}, \quad \Omega_{\rho, \tau, \nu}^{\infty} \approx \Lambda_{\rho, \tau, \nu}, \\
& b \Omega_{\rho, \tau, \nu} \approx b \Lambda_{\rho, \tau, \nu},  \tag{5.19}\\
& Z_{p, \rho}^{\infty}=\left\{k \in \mathbb{C}^{3}:(k, p) \in \Omega_{\nu} \backslash \bar{\Omega}_{\rho}\right\} \approx \mathcal{D}_{\rho /|p|} \text { for fixed } p,
\end{align*}
$$

where $\rho>0,0<\tau<1, \nu \in \mathbb{S}^{2}$ (and where we use the definitions (5.3), (5.6), (5.13)).
In [19] $\lambda, p$ of (5.11) were used as coordinates on $\Omega$. In the present work we use them also as coordinates on $\Omega \backslash \Omega_{\rho}$ (or more precisely on $\Omega_{\nu} \backslash \Omega_{\rho}$ ).

## 6. An integral equation of [20] and some related formulas

In the main considerations of [20] it is assumed that $d=3$ and the basic assumption on $v$ consists in the following condition on its Fourier transform:

$$
\begin{equation*}
\hat{v} \in L_{\mu}^{\infty}\left(\mathbb{R}^{3}\right) \cap \mathcal{C}\left(\mathbb{R}^{3}\right) \text { for some real } \mu \geq 2 \tag{6.1}
\end{equation*}
$$

where $\hat{v}$ is defined by (3.12) (for $d=3$ ),

$$
\begin{align*}
& L_{\mu}^{\infty}\left(\mathbb{R}^{d}\right)=\left\{u \in L^{\infty}\left(\mathbb{R}^{d}\right): \quad\|u\|_{\mu}<+\infty\right\} \\
& \|u\|_{\mu}=\text { ess } \sup _{p \in \mathbb{R}^{d}}(1+|p|)^{\mu}|u(p)|, \quad \mu>0 \tag{6.2}
\end{align*}
$$

and $\mathcal{C}$ denotes the space of continuous functions.
Note that

$$
\begin{align*}
& v \in W^{m, 1}\left(\mathbb{R}^{d}\right) \Longrightarrow \hat{v} \in L_{\mu}^{\infty}\left(\mathbb{R}^{d}\right) \cap \mathcal{C}\left(\mathbb{R}^{d}\right)  \tag{6.3}\\
& \|\hat{v}\|_{\mu} \leq c_{4}(m, d)\|v\|_{m, 1} \text { for } \mu=m
\end{align*}
$$

where $W^{m, 1}, L_{\mu}^{\infty}$ are the spaces of (2.2), (6.2).
Let

$$
\begin{equation*}
H(\lambda, p)=H(k(\lambda, p), p), \quad(\lambda, p) \in(\mathbb{C} \backslash 0) \times\left(\mathbb{R}^{3} \backslash \mathcal{L}_{\mu}\right) \tag{6.4}
\end{equation*}
$$

where $H$ is the function of (3.13), $\lambda, p$ are the coordinates of Subsection 5.2 under assumption (5.9a).

Let

$$
\begin{align*}
& L_{\mu}^{\infty}\left(\Lambda_{\rho, \tau, \nu}\right)=\left\{U \in L^{\infty}\left(\Lambda_{\rho, \tau, \nu}\right):\|U \mid\| \|_{\rho, \tau, \mu}<\infty\right\} \\
& \left\|\left|\|U \mid\| \|_{\rho, \tau, \mu}=\text { ess } \sup _{(\lambda, p) \in \Lambda_{\rho, \tau, \nu}}(1+|p|)^{\mu}\right| U(\lambda, p) \mid, \quad \mu>0\right. \tag{6.5}
\end{align*}
$$

where $\Lambda_{\rho, \tau, \nu}$ is defined in (5.13), $\left.\rho>0, \tau \in\right] 0,1\left[, \nu \in \mathbb{S}^{2}, \mu>0\right.$.
Let $v$ satisfy (6.1) and $\|\hat{v}\|_{\mu} \leq C$. Let

$$
\begin{equation*}
\eta(C, \rho, \mu) \stackrel{\text { def }}{=} a(\mu) C(\ln \rho)^{2} \rho^{-1}<1, \ln \rho \geq 2 \tag{6.6}
\end{equation*}
$$

where $a(\mu)$ is the constant $c_{2}(\mu)$ of [20]. Let $H(\lambda, p)$ be defined by (6.4) and be considered as a function on $\Lambda_{\rho, \tau, \nu}$ of (5.13). Then (see Section 4 of [20]):

$$
\begin{equation*}
\left.H=H^{0}+M_{\rho, \tau}(H)+Q_{\rho, \tau}, \quad \tau \in\right] 0,1[ \tag{6.7}
\end{equation*}
$$

where

$$
\begin{align*}
& H^{0}(\lambda, p)=\frac{1}{2 \pi i} \int_{\mathcal{T}_{\rho \backslash|p|}^{+}} H(\zeta, p) \frac{d \zeta}{\zeta-\lambda},(\lambda, p) \in \Lambda_{\rho, \tau, \nu}^{+}  \tag{6.8a}\\
& H^{0}(\lambda, p)=-\frac{1}{2 \pi i} \int_{\mathcal{T}_{\rho /|p|}^{-}} H(\zeta, p) \frac{\lambda d \zeta}{\zeta(\zeta-\lambda)},(\lambda, p) \in \Lambda_{\rho, \tau, \nu}^{-} \tag{6.8b}
\end{align*}
$$

where $\Lambda_{\rho, \tau, \nu}^{ \pm}, \mathcal{T}_{r}^{ \pm}$are defined in (5.17), (5.18) (and where the integrals along $\mathcal{T}_{r}^{ \pm}$are taken in the counter-clock wise direction);

$$
\begin{align*}
& M_{\rho, \tau}(U)(\lambda, p)=M_{\rho, \tau}^{+}(U)(\lambda, p)= \\
& -\frac{1}{\pi} \iint_{\mathcal{D}_{\rho /|p|}^{+}}(U, U)_{\rho, \tau}(\zeta, p) \frac{d \operatorname{Re} \zeta d \operatorname{Im} \zeta}{\zeta-\lambda},(\lambda, p) \in \Lambda_{\rho, \tau, \nu}^{+}  \tag{6.9a}\\
& M_{\rho, \tau}(U)(\lambda, p)=M_{\rho, \tau}^{-}(U)(\lambda, p)= \\
& -\frac{1}{\pi} \int_{\mathcal{D}_{\rho /|p|}^{-}}(U, U)_{\rho, \tau}(\zeta, p) \frac{\lambda d \operatorname{Re} \zeta d I m \zeta}{\zeta(\zeta-\lambda)},(\lambda, p) \in \Lambda_{\rho, \tau, \nu}^{-}  \tag{6.9b}\\
& \left(U_{1}, U_{2}\right)_{\rho, \tau}(\zeta, p)=\left\{\chi_{2 \tau \rho} U_{1}^{\prime}, \chi \chi_{2 \tau \rho} U_{2}^{\prime}\right\}(\zeta, p),(\zeta, p) \in \Lambda_{\rho, \tau, \nu} \\
& \chi_{2 \tau \rho} U_{j}^{\prime}(k, p)=U_{j}(\lambda(k, p), p), \quad(k, p) \in \Omega_{\rho, \tau, \nu}^{\infty},  \tag{6.10}\\
& \chi_{2 \tau \rho} U_{j}^{\prime}(k, p)=0, \quad|p| \geq 2 \tau \rho, \quad j=1,2
\end{align*}
$$

where $U, U_{1}, U_{2}$ are test functions on $\Lambda_{\rho, \tau, \nu}, \Omega_{\rho, \tau, \nu}^{\infty}$ is defined in (5.6), $\lambda(k, p)$ is defined in (5.11a), $\{\cdot, \cdot\}$ is defined by the formula

$$
\begin{align*}
& \left\{F_{1}, F_{2}\right\}(\lambda, p)=-\frac{\pi}{4} \int_{-\pi}^{\pi}\left(\frac{|p|}{2} \frac{|\lambda|^{2}-1}{\bar{\lambda}|\lambda|}(\cos \varphi-1)-\frac{|p|}{\bar{\lambda}} \sin \varphi\right) \times  \tag{6.11}\\
& F_{1}(k(\lambda, p),-\xi(\lambda, p, \varphi)) F_{2}(k(\lambda, p)+\xi(\lambda, p, \varphi), p+\xi(\lambda, p, \varphi)) d \varphi
\end{align*}
$$

for $(\lambda, p) \in \Lambda_{\rho, \nu}$, where $F_{1}, F_{2}$ are test functions on $\Omega \backslash \bar{\Omega}_{\rho}, k(\lambda, p)$ is defined in (5.11b), $\Lambda_{\rho, \nu}$ is defined in (5.13),

$$
\begin{align*}
& \xi(\lambda, p, \varphi)=\operatorname{Re} k(\lambda, p)(\cos \varphi-1)+k^{\perp}(\lambda, p) \sin \varphi  \tag{6.12}\\
& k^{\perp}(\lambda, p)=\frac{\operatorname{Im} k(\lambda, p) \times \operatorname{Re} k(\lambda, p)}{|\operatorname{Im} k(\lambda, p)|} \tag{6.13}
\end{align*}
$$

where $\times$ in (6.13) denotes vector product;

$$
\begin{align*}
& H, H^{0}, Q_{\rho, \tau} \in L_{\mu}^{\infty}\left(\Lambda_{\rho, \tau, \nu}\right)  \tag{6.14}\\
& \left\|\|H\|_{\rho, \tau, \mu_{0}} \leq \frac{C}{1-\eta(C, \rho, \mu)}\right.  \tag{6.15a}\\
& \left\|\left\|H^{0}\right\|\right\|_{\rho, \tau, \mu_{0}} \leq \frac{C}{1-\eta(C, \rho, \mu)}\left(1+\frac{c_{5}\left(\mu_{0}\right) C}{1-\eta(C, \rho, \mu)}\right)  \tag{6.15b}\\
& \left\|\mid Q_{\rho, \tau}\right\| \|_{\rho, \tau, \mu_{0}} \leq \frac{3 c_{5}\left(\mu_{0}\right) C^{2}}{(1-\eta(C, \rho, \mu))^{2}(1+2 \tau \rho)^{\mu-\mu_{0}}} \tag{6.15c}
\end{align*}
$$

where $2 \leq \mu_{0} \leq \mu, c_{5}$ is the constant $b_{4}$ of [20], $\eta(C, \rho, \mu)$ is defined by (6.6).

Following [20] we consider (6.7) as an approxiate integral equation for finding $H$ on $\Lambda_{\rho, \tau, \nu}$ from $H^{0}$ on $\Lambda_{\rho, \tau, \nu}$ with unknown remainder $Q_{\rho, \tau}$.

Note also that if $\hat{v}$ satisfies (6.1), then (see [19], [20])

$$
\begin{align*}
& H(\lambda, p) \rightarrow \hat{v}(p) \text { as } \lambda \rightarrow 0  \tag{6.16}\\
& H(\lambda, p) \rightarrow \hat{v}(p) \text { as } \lambda \rightarrow \infty
\end{align*}
$$

where $p \in \mathcal{B}_{2 \tau \rho} \backslash \mathcal{L}_{\nu}, H$ is defined by (6.4) and is considered as a function on $\Lambda_{\rho, \tau, \nu}, \rho>0$, $0<\tau<1, \nu \in \mathbb{S}^{2}$.

## 7. Proof of Theorem 2.2 for $d=3$

Lemma 7.1. Let $\hat{v}_{i}$ satisfy (6.1) and $\left\|\hat{v}_{i}\right\|_{\mu}<C$, where $i=1,2$. Let

$$
\begin{equation*}
0<\tau \leq \tau_{1}\left(\mu, \mu_{0}, C, \delta\right), \rho \geq \rho_{1}\left(\mu, \mu_{0}, C, \delta\right) \tag{7.1}
\end{equation*}
$$

where $\tau_{1}, \rho_{1}$ are the constants of Section 4 of [20] and where $\delta=1 / 2,2 \leq \mu_{0}<\mu$. Then

$$
\begin{equation*}
\left\|\left\|H_{2}-H_{1}\right\|\right\|_{\rho, \tau, \mu_{0}} \leq 2\left(\left\|\left|H_{2}^{0}-H_{1}^{0}\| \|_{\rho, \tau, \mu_{0}}+\left\|\left|Q_{\rho, \tau}^{2}-Q_{\rho, \tau}^{1}\right|\right\| \|_{\rho, \tau, \mu_{0}}\right),\right.\right. \tag{7.2}
\end{equation*}
$$

where $H_{i}, H_{i}^{0}, Q_{\rho, \tau}^{i}$ are the functions of (6.4), (6.7), (6.8), (6.14), (6.15) for $v=v_{i}, i=1,2$. In addition,

$$
\begin{equation*}
\left\|\left|\left|Q_{\rho, \tau}^{2}-Q_{\rho, \tau}^{1}\right| \|_{\rho, \tau, \mu_{0}} \leq \frac{24 c_{5}\left(\mu_{0}\right) C^{2}}{(1+2 \tau \rho)^{\mu-\mu_{0}}}\right.\right. \tag{7.3}
\end{equation*}
$$

In connection with (7.1) we remind that $\left.\tau_{1} \in\right] 0,1\left[\right.$ is sufficiently small and $\rho_{1}$ is sufficiently great, see [20].

Lemma 7.1 follows from estimates mentioned as estimates (6.14), (6.15) of the present paper (see estimates (3.3), (4.20), (4.22) of [20]) and from Lemmas 4.4, 4.5 and estimates (4.36) of [20].

Lemma 7.2. Let $\hat{v}_{i}$ satisfy (6.1) and $\left\|\hat{v}_{i}\right\|_{\mu}<C$, where $i=1,2$. Let

$$
\begin{equation*}
\eta\left(C, \rho_{0}, \mu\right) \leq 1 / 2, \quad \ln \rho_{0} \geq 2 \tag{7.4}
\end{equation*}
$$

where $\eta$ is defined in (6.6). Let

$$
0<\tau_{0}<1, \quad 2 \leq \mu_{0}<\mu, \quad \rho=2 \rho_{0}, \quad \tau=\tau_{0} / 2
$$

Then

$$
\begin{equation*}
\left.\left\|\left|H_{2}^{0}-H_{1}^{0}\right|\right\|\right|_{\rho, \tau, \mu_{0}} \leq\left(c_{6}+4 c_{7}\left(\mu_{0}, \tau_{0}, \rho_{0}\right) C\right) \mid\left\|\chi_{\rho_{0}, \tau_{0}, \rho, \tau}\left(H_{2}-H_{1}\right)\right\| \|_{\rho_{0}, \tau_{0}, \mu_{0}}, \tag{7.5}
\end{equation*}
$$

where $H_{i}, H_{i}^{0}$ are the functions of (6.4), (6.8) for $v=v_{i}, i=1,2, \chi_{\rho_{0}, \tau_{0}, \rho, \tau}$ is the characteristic function of $\Lambda_{\rho_{0}, \tau_{0}, \nu} \backslash \Lambda_{\rho, \tau, \nu}, c_{6}$ is defined by (8.9), $c_{7}$ is the constant $c_{8}$ of [20] (that is $c_{7}(\mu, \tau, \rho)=3 b_{1}(\mu) \tau^{2}+4 b_{2}(\mu) \rho^{-1}+4 b_{3}(\mu) \tau$, where $b_{1}, b_{2}, b_{3}$ are the constants of [20]).

Note that
$\Lambda_{\rho, \tau, \nu} \subset \Lambda_{\rho_{0}, \tau_{0}, \nu}$ under the assumptions that $\rho=2 \rho_{0}, \tau=\tau_{0} / 2,0<\tau_{0}<1, \rho>0$.
Lemma 7.2 is proved in Section 8.
Lemma 7.3. Let the assumptions of Theorem 2.1 hold (for $d=3$ ). Let

$$
\begin{equation*}
\rho_{0} \geq r_{1}(N, D, m, \sigma) \text { for some } \sigma>1 \tag{7.7}
\end{equation*}
$$

where $r_{1}$ is the number of (3.9). Let $0<\tau_{0}<1,0<\mu_{0}, \rho=2 \rho_{0}, \tau=\tau_{0} / 2$. Then

$$
\begin{equation*}
\left\|\left\|\chi_{\rho_{0}, \tau_{0}, \rho, \tau}\left(H_{2}-H_{1}\right)\right\|\right\|_{\rho_{0}, \tau_{0}, \mu_{0}} \leq c_{8} \sigma^{2} e^{2 \rho L}\left\|\Phi_{2}-\Phi_{1}\right\|(1+\rho)^{\mu_{0}} \tag{7.8}
\end{equation*}
$$

where

$$
\begin{equation*}
c_{8}=(2 \pi)^{-d} \int_{\partial D} d x, \quad L=\max _{x \in \partial D}|x|, \tag{7.9}
\end{equation*}
$$

$\left\|\Phi_{2}-\Phi_{1}\right\|$ is defined according to (2.4), $\chi_{\rho_{0}, \tau_{0}, \rho, \tau}, H_{1}, H_{2}$ are the same that in (7.5).
Lemma 7.3 is proved in Section 8.
Lemma 7.4. Let the assumptions of Theorem 2.1 hold (for $d=3$ ). Let

$$
\begin{equation*}
0<\tau \leq \tau_{2}\left(m, \mu_{0}, N\right), \quad \rho \geq \rho_{2}\left(m, \mu_{0}, N, D, \sigma\right), \tag{7.10}
\end{equation*}
$$

where $2 \leq \mu_{0}<m, \sigma>1$, and $\tau_{2}, \rho_{2}$ are constants such that (7.10) implies that

$$
\begin{align*}
& \tau \leq \tau_{1}\left(m, \mu_{0}, c_{4}(m, 3) N, 1 / 2\right), \quad \rho \geq \rho_{1}\left(m, \mu_{0}, c_{4}(m, 3) N, 1 / 2\right)  \tag{7.11a}\\
& \tau<1 / 2, \eta\left(c_{4}(m, 3) N, \rho / 2, m\right) \leq 1 / 2, \ln (\rho / 2) \geq 2  \tag{7.11b}\\
& \rho / 2 \geq r_{1}(N, D, m, \sigma) \tag{7.11c}
\end{align*}
$$

where $\tau_{1}, \rho_{1}, \eta, r_{1}$ are the same that in (7.1), (7.4), (7.7), $c_{4}$ is the constant of (6.3). Then

$$
\begin{equation*}
\left\|H_{2}-H_{1}\right\|\left\|_{\rho, \tau, \mu_{0}} \leq c_{9}\left(N, D, m, \mu_{0}, \sigma, \tau\right) e^{2 L \rho} \rho^{\mu_{0}}\right\| \Phi_{2}-\Phi_{1} \|+c_{10}\left(N, m, \mu_{0}, \tau\right) \rho^{-\left(m-\mu_{0}\right)} \tag{7.12}
\end{equation*}
$$

where $c_{9}, c_{10}$ are some constants which can be given explicitly.
Lemma 7.4 follows from formula (6.3) and Lemmas 7.1, 7.2, 7.3.
The final part of the proof of Theorem 2.2 for $d=3$ consists of the following. Under the assumptions of Lemma 7.4 for $\mu_{0}=2$, we have that

$$
\begin{align*}
& \left\|v_{1}-v_{2}\right\|_{L^{\infty}(D)} \leq\left\|\left|\left|H_{2}-H_{1}\right| \|_{\rho, \tau, 2} \int_{|p| \leq 2 \rho \tau} \frac{d p}{(1+|p|)^{2}}+\right.\right. \\
& 2 c_{4}(m, 3) N \int_{|p| \geq 2 \rho \tau} \frac{d p}{(1+|p|)^{m}} \leq  \tag{7.13}\\
& 8 \pi \rho \tau\left|\left\|H_{2}-H_{1} \mid\right\|_{\rho, \tau, 2}+\frac{8 \pi c_{4}(m, 3) N}{(m-3)(2 \tau)^{m-3}} \frac{1}{\rho^{m-3}} \leq\right. \\
& c_{11}(N, D, m, \sigma, \tau) e^{2 L \rho} \rho^{3}\left\|\Phi_{2}-\Phi_{1}\right\|+c_{12}(N, m, \tau) \rho^{-(m-3)}
\end{align*}
$$

where $c_{11}, c_{12}$ are related in a simple way with $c_{9}, c_{10}$ for $\mu_{0}=2$. To obtain (7.13) we used also (6.3), (6.5), (6.16) and the inverse Fourier transform formula

$$
\begin{equation*}
v(x)=\int_{\mathbb{R}^{3}} e^{-i p x} \hat{v}(p) d p, \quad x \in \mathbb{R}^{3} \tag{7.14}
\end{equation*}
$$

Let now

$$
\begin{equation*}
\alpha \in] 0,1\left[, \beta=\frac{1-\alpha}{2 L}, \delta=\left\|\Phi_{1}-\Phi_{2}\right\|, \rho=\beta \ln \left(1+\delta^{-1}\right)\right. \tag{7.15}
\end{equation*}
$$

where $\delta$ is so small that $\rho \geq \rho_{2}(m, 2, N, D, \sigma)$, where $\rho_{2}$ is the constant of (7.10). Then, due to (7.13),

$$
\begin{align*}
& \left\|v_{1}-v_{2}\right\|_{L^{\infty}(D)} \leq c_{11}(N, D, m, \sigma, \tau)\left(1+\delta^{-1}\right)^{2 L \beta}\left(\beta \ln \left(1+\delta^{-1}\right)\right)^{3} \delta+ \\
& c_{12}(N, D, m, \tau)\left(\beta \ln \left(1+\delta^{-1}\right)\right)^{-(m-3)}= \\
& c_{11}(N, D, m, \sigma, \tau) \beta^{3}(1+\delta)^{1-\alpha} \delta^{\alpha}\left(\ln \left(1+\delta^{-1}\right)\right)^{3}+  \tag{7.16}\\
& c_{12}(N, D, m, \tau) \beta^{-(m-3)}\left(\ln \left(1+\delta^{-1}\right)\right)^{-(m-3)},
\end{align*}
$$

where $\sigma, \tau$ are the same that in (7.10) for $\mu_{0}=2$ and where $\alpha, \beta$ and $\delta$ are the same that in (7.15).

Using (7.16) we obtain that

$$
\begin{equation*}
\left\|v_{1}-v_{2}\right\|_{L^{\infty}(D)} \leq c_{13}(N, D, m)\left(\ln \left(1+\left\|\Phi_{1}-\Phi_{2}\right\|^{-1}\right)\right)^{-(m-3)} \tag{7.17}
\end{equation*}
$$

for $\delta=\left\|\Phi_{1}-\Phi_{2}\right\| \leq \delta_{0}(N, D, m)$, where $\delta_{0}$ is sufficiently small positive constant. Estimate (7.17) in the general case (with modified $c_{13}$ ) follows from (7.17) for $\delta=\left\|\Phi_{1}-\Phi_{2}\right\| \leq \delta_{0}(N, D, m)$ and the property that $\left\|v_{j}\right\|_{L^{\infty}(D)} \leq c_{14}(m) N($ for $d=3)$.

Thus, Theorem 2.2 for $d=3$ is proved.

## 8. Proofs of Lemmas 7.2 and 7.3

Proof of Lemma 7.2. Using the maximum principle for holomorphic functions it is sufficient to prove that

$$
\begin{align*}
& \sup _{(\lambda, p) \in b \Lambda_{\rho, \tau, \nu}^{ \pm}}(1+|p|)^{\mu_{0}}\left|H_{2}^{0}(\lambda(1 \mp 0), p)-H_{1}^{0}(\lambda(1 \mp 0), p)\right| \leq  \tag{8.1}\\
& \left(c_{6}+4 c_{7}\left(\mu_{0}, \tau_{0}, \rho_{0}\right) C\right)\left|\| \chi_{\rho_{0}, \tau_{0}, \rho, \tau}\left(H_{2}-H_{1}\right)\right|\left|\left.\right|_{\rho_{0}, \tau_{0}, \mu_{0}},\right.
\end{align*}
$$

where $b \Lambda_{\rho, \tau, \nu}^{+}, b \Lambda_{\rho, \tau, \nu}^{-}$are defined in (5.17) (and where $H_{i}^{0}(\lambda(1-0), p), i=1,2$, are considered for $(\lambda, p) \in b \Lambda_{\rho, \tau, \nu}^{+}, H_{i}^{0}(\lambda(1+0), p), i=1,2$, are considered for $\left.(\lambda, p) \in b \Lambda_{\rho, \tau, \nu}^{-}\right)$.

Using (6.8) and the Sohotsky-Plemelj formula, we have that

$$
\begin{align*}
& H^{0}(\lambda(1-0), p)=\frac{1}{2 \pi i} \int_{\mathcal{T}_{\rho /|p|}^{+}} H(\zeta, p) \frac{d \zeta}{\zeta-\lambda(1+0)}+H(\lambda, p),(\lambda, p) \in b \Lambda_{\rho, \tau, \nu}^{+},  \tag{8.2a}\\
& H^{0}(\lambda(1+0), p)=-\frac{1}{2 \pi i} \int_{\mathcal{T}_{\rho| | p \mid}^{-}} H(\zeta, p) \frac{\lambda d \zeta}{\zeta(\zeta-\lambda(1-0))}+H(\lambda, p),(\lambda, p) \in b \Lambda_{\rho, \tau, \nu}^{-} \tag{8.2b}
\end{align*}
$$

In addition, using the Cauchy-Green formula we have that

$$
\begin{align*}
& H(\lambda, p)=-\frac{1}{2 \pi i} \int_{\mathcal{T}_{\rho /|p|}^{+}} H(\zeta, p) \frac{d \zeta}{\zeta-\lambda(1+0)}+\frac{1}{2 \pi i} \int_{\mathcal{T}_{\rho /|p|}^{+}} H(\zeta, p) \frac{d \zeta}{\zeta-\lambda}-  \tag{8.3a}\\
& \frac{1}{\pi} \int_{\mathcal{D}_{\rho_{0} /|p|}^{+}} \int_{\mathcal{D}_{\rho /|p|}^{+}} \frac{\partial H(\zeta, p)}{\partial \bar{\zeta}} \frac{d \operatorname{Re} \zeta d \operatorname{Im} \zeta}{\zeta-\lambda},(\lambda, p) \in b \Lambda_{\rho, \tau, \nu}^{+}, \\
& H(\lambda, p)=\frac{1}{2 \pi i} \int_{\mathcal{T}_{\rho /|p|}^{-}}^{-} H(\zeta, p) \frac{\lambda d \zeta}{\zeta(\zeta-\lambda(1-0))}-\frac{1}{2 \pi i} \int_{\mathcal{T}_{\rho /|p|}^{-}} H(\zeta, p) \frac{\lambda d \zeta}{\zeta(\zeta-\lambda)}-  \tag{8.3b}\\
& \frac{1}{\pi} \int_{\mathcal{D}_{\rho_{0} /|p|}^{-} \backslash \mathcal{D}_{\rho /|p|}^{-}} \int \frac{\partial H(\zeta, p)}{\partial \bar{\zeta}} \frac{\lambda d \operatorname{Re} \zeta d I m \zeta}{\zeta(\zeta-\lambda)},(\lambda, p) \in b \Lambda_{\rho, \tau, \nu}^{-}
\end{align*}
$$

(where the integrals along $\mathcal{T}_{r}^{ \pm}$are taken in the counter- clockwise direction). In addition (see formulas (3.22), (3.23), (4.8), (4.14) of [20]),

$$
\begin{equation*}
\frac{\partial H(\zeta, p)}{\partial \bar{\zeta}}=(H, H)_{\rho_{0}, \tau_{0}}(\zeta, p), \quad(\zeta, p) \in \Lambda_{\rho_{0}, \tau_{0}, \nu} \tag{8.4}
\end{equation*}
$$

where $(\cdot, \cdot)_{\rho, \tau}$ is defined by (6.10).
Using (8.2), (8.3) we obtain that

$$
\begin{align*}
& H^{0}(\lambda(1-0), p)=\frac{1}{2 \pi i} \int_{\mathcal{T}_{\rho_{0} /|p|}^{+}} H(\zeta, p) \frac{d \zeta}{\zeta-\lambda}-  \tag{8.5a}\\
& \frac{1}{\pi} \int_{\mathcal{D}_{\rho_{0} /|p|}^{+} \backslash \mathcal{D}_{\rho /|p|}^{+}} \frac{\partial H(\zeta, p)}{\partial \bar{\zeta}} \frac{d \operatorname{Re} \zeta d I m \zeta}{\zeta-\lambda},(\lambda, p) \in b \Lambda_{\rho, \tau, \nu}^{+} \\
& H^{0}(\lambda(1+0), p)=-\frac{1}{2 \pi i} \int H(\zeta, p) \frac{\lambda d \zeta}{\zeta(\zeta-\lambda)}-  \tag{8.5b}\\
& \frac{1}{\pi} \int_{\mathcal{D}_{\rho_{0} /||p|}^{-} \backslash \mathcal{D}_{\rho /|p|}^{-}} \int_{\mathcal{T}_{\rho_{0} /|p|}^{-}} \frac{\partial H(\zeta, p)}{\partial \bar{\zeta}} \frac{\lambda d \operatorname{Re} \zeta d I m \zeta}{\zeta(\zeta-\lambda)},(\lambda, p) \in b \Lambda_{\rho, \tau, \nu}^{-} .
\end{align*}
$$

Using (8.5), (8.4) for $H^{0}=H_{i}^{0}, H=H_{i}, i=1,2$, we obtain that:

$$
\begin{align*}
& \left(H_{2}^{0}-H_{1}^{0}\right)(\lambda(1-0), p)=A^{+}(\lambda, p)+B^{+}(\lambda, p)  \tag{8.6a}\\
& A^{+}(\lambda, p)=\frac{1}{2 \pi i} \int_{\mathcal{T}_{\rho_{0} /|p|}^{+}}\left(H_{2}-H_{1}\right)(\zeta, p) \frac{d \zeta}{\zeta-\lambda}, \\
& B^{+}(\lambda, p)=-\frac{1}{\pi} \int_{\mathcal{D}_{\rho_{0} /|p|}^{+} \backslash \mathcal{D}_{\rho /|p|}^{+}}\left(\left(H_{2}-H_{1}, H_{2}\right)_{\rho_{0}, \tau_{0}}+\left(H_{1}, H_{2}-H_{1}\right)_{\rho_{0}, \tau_{0}}\right)(\zeta, p) \frac{d \operatorname{Re} \zeta d \operatorname{Im} \zeta}{\zeta-\lambda}
\end{align*}
$$

where $(\lambda, p) \in b \Lambda_{\rho, \tau, \nu}^{+}$;

$$
\begin{align*}
& \left(H_{2}^{0}-H_{1}^{0}\right)(\lambda(1+0), p)=A^{-}(\lambda, p)+B^{-}(\lambda, p),  \tag{8.6b}\\
& A^{-}(\lambda, p)=-\frac{1}{2 \pi i} \int_{\mathcal{T}_{\rho_{0} /|p|}}\left(H_{2}-H_{1}\right)(\zeta, p) \frac{\lambda d \zeta}{\zeta(\zeta-\lambda)}, \\
& B^{-}(\lambda, p)=-\frac{1}{\pi} \int_{\mathcal{D}_{\rho_{0} /|p|}^{-}} \int_{\mathcal{D}_{\rho /|p|}^{-}}\left(\left(H_{2}-H_{1}, H_{2}\right)_{\rho_{0}, \tau_{0}}+\left(H_{1}, H_{2}-H_{1}\right)_{\rho_{0}, \tau_{0}}\right)(\zeta, p) \frac{\lambda d \operatorname{Re} \zeta d \operatorname{Im} \zeta}{\zeta(\zeta-\lambda)},
\end{align*}
$$

where $(\lambda, p) \in b \Lambda_{\rho, \tau, \nu}^{-}$.
Estimates (8.1) follow from formulas (8.6) and from the estimates

$$
\begin{align*}
& \left|A^{ \pm}(\lambda, p)\right| \leq c_{6}(1+|p|)^{-\mu_{0}} \Delta, \quad(\lambda, p) \in b \Lambda_{\rho, \tau, \nu}^{ \pm},  \tag{8.7}\\
& \left|B^{ \pm}(\lambda, p)\right| \leq 4 c_{7}\left(\mu_{0}, \tau_{0}, \rho_{0}\right) C \Delta, \quad(\lambda, p) \in b \Lambda_{\rho, \tau, \nu}^{ \pm} \tag{8.8}
\end{align*}
$$

where

$$
\begin{gather*}
c_{6}=\sup _{r \in] 1 / 2,+\infty} \frac{q(r)}{q(r)-q(2 r)},  \tag{8.9}\\
q(r)=2 r\left(1-\left(1-\frac{1}{4 r^{2}}\right)^{1 / 2}\right) \\
\Delta=\| \| \chi_{\rho_{0}, \tau_{0}, \rho, \tau}\left(H_{2}-H_{1}\right) \mid \|_{\rho_{0}, \tau_{0}, \mu_{0}} . \tag{8.10}
\end{gather*}
$$

Note that

$$
\begin{equation*}
0<c_{6} \leq(2 \sqrt{3}-3)^{-1} \tag{8.11}
\end{equation*}
$$

where $c_{6}$ is defined by (8.9). Estimate (8.11) follows from the formulas

$$
\begin{align*}
& c_{6}=\frac{1}{1-2 \sigma}, \quad \sigma=\sup _{\tau \in] 0,1[ } \frac{1-(1-(1 / 4) \tau)^{1 / 2}}{1-(1-\tau)^{1 / 2}},  \tag{8.12}\\
& (1-\tau)^{1 / 2} \leq 1-(1 / 2) \tau, \quad 1-(1 / 4) \tau \geq a(1-(1 / 4) \tau)+1-a,  \tag{8.13}\\
& a=2(2-\sqrt{3}), \tau \in] 0,1[.
\end{align*}
$$

Estimates (8.7) follow from formula (7.6), the properties that

$$
\begin{align*}
& H_{i} \in \mathcal{C}\left(\Lambda_{\rho_{0}, \tau_{0}, \nu} \cup b \Lambda_{\rho_{0}, \tau_{0}, \nu}\right) \\
& \left|H_{i}(\lambda, p)\right| \leq(1+|p|)^{-\mu}| | H_{i} \mid \|_{\rho_{0}, \tau_{0}, \mu}, \quad(\lambda, p) \in \Lambda_{\rho_{0}, \tau_{0}, \nu} \cup b \Lambda_{\rho_{0}, \tau_{0}, \nu}, i=1,2, \tag{8.14}
\end{align*}
$$

(see formulas (3.2), (3.3) of [20] for $H$ and formulas (5.3), (5.6), (5.13), (5.19), (6.4), (6.5) of the present paper for $\Omega_{\rho, \tau, \nu}^{\infty}, \Lambda_{\rho, \tau, \nu}, H$ and $\left.L_{\mu}^{\infty}\left(\Lambda_{\rho, \tau, \nu}\right)\right)$ and from the formulas

$$
\begin{equation*}
\frac{1}{2 \pi} \int_{\mathcal{T}_{r}^{-}} \frac{|\lambda||d \zeta|}{|\zeta||\zeta-\lambda|}=\frac{1}{2 \pi} \int_{\mathcal{T}_{r}^{+}} \frac{\left|d \zeta^{\prime}\right|}{\left|\zeta^{\prime}-\lambda^{\prime}\right|} \leq \frac{q(r)}{q(r)-q(2 r)} \tag{8.15}
\end{equation*}
$$

$\lambda \in \mathcal{T}_{2 r}^{-}, \lambda^{\prime} \in \mathcal{T}_{2 r}^{+}, r>1 / 2$. In turn, formulas (8.15) follow from the property that $z^{-1} \in \mathcal{T}_{r}^{+}$if $z \in \mathcal{T}_{r}^{-}$and from the formula that $q(r)$ is the radius of $\mathcal{T}_{r}^{+}$, where $r>1 / 2$.

Estimates (8.8) follow from the proof of estimates (7.8), (7.9) of [20] and from the formulas (6.15a) (for $\rho=\rho_{0}, \tau=\tau_{0}$ ), (7.4), (7.6), (8.14) of the present paper.

Lemma 7.2 is proved.
Proof of Lemma 7.3. Using formulas (6.4), (6.5), the formulas

$$
\Lambda_{\rho, \tau, \nu} \approx \Omega_{\rho, \tau, \nu}^{\infty}=\Omega_{\rho, \tau}^{\infty} \cap \Omega_{\nu}, \quad \Omega_{\rho, \tau}^{\infty} \approx \Theta_{\rho, \tau}^{\infty}
$$

(see (5.19), (5.6), (5.4)), and formulas (5.2), (3.13), we have that

$$
\begin{equation*}
\Delta \leq \sup _{\substack{(k, l) \in \bar{\Theta}_{\rho} \backslash \Theta_{\rho_{0}},|k-l| \leq 2 \tau \rho}}\left|h_{2}(k, l)-h_{1}(k, l)\right|, \tag{8.16}
\end{equation*}
$$

where $\Delta$ is defined by (8.10), $h_{1}, h_{2}$ are the functions of Section 4.
Estimate (7.8) follows from formulas (8.16), (4.1), (3.7), (3.9), the property that $|k-l| \leq \rho$ in (8.16) (since $\tau<1 / 2$ due to the assumptions of Lemma 7.3) and the property that $\left|e^{i k x}\right| \leq e^{\rho L},\left|e^{i l x}\right| \leq e^{\rho L}$ for $k, l \in \bar{\Theta}_{\rho}, x \in \partial D$.

Lemma 7.3 is proved.

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