# Fundamental solutions of the acoustic and diffusion equations in nonhomogeneous medium 

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#### Abstract

A fundamental solution of the acoustical equation with a variable refraction coefficient is constructed. The solution satisfies the limiting absorption and radiation conditions. The optimal high frequency estimate is proved for square means of the solution. The source function for the diffusion equation is a by-product of this construction.


## 1. Introduction

Consider the generalized Helmholtz equation

$$
\begin{equation*}
\left(\Delta+\omega^{2} \mathbf{n}^{2}+\mathbf{m}\right) u=0 \tag{1}
\end{equation*}
$$

in a Euclidean space $X$ with time frequency $\dot{u}$ and real variable refraction coefficients $\mathbf{n}$ and "mass" $\mathbf{m}$. The acoustical (wave) equation in frequency domain is the particular case $(\mathbf{m}=0)$. A fundamental solution for (1) is a function $S(y, x ; \omega)$ in $X \times X$ that satisfies the equation

$$
\left(\Delta_{x}+\omega^{2} \mathbf{n}^{2}(x)+\mathbf{m}(x)\right) S(y . x ; \omega)=\delta(x-y) .
$$

We say that the fundamental solution satisfies the limiting absorption condition, if it admits an analytic continuation to the halfplane $\mathbf{C}_{+}\left(\omega_{0}\right) \doteq\left\{\omega: \operatorname{Im} \omega>0,|\dot{\omega}|>\omega_{0}\right\}$ for some $\omega_{0}>0$ that tends to zero, as $|\omega| \rightarrow \infty$. We call such a fundamental solution a source function of (1). The uniqueness of the source function is easy to check (see below). The source function describes a divergent time harmonic wave with the phase function $\phi=-\omega t$. We here state the existence of a source function for a smooth medium in the plane $X$, which is homogeneous outside a compact set. We

[^0]also give the estimate (Section 3) which shows the rate of decrease of the weighted square means of the source function, as $|\omega| \rightarrow \infty$. This information is applied to the analysis of near scattering fields, see [8].

The diffusion equation in the plane can be reduced to (1) and we show that the decreasing fundamental solution can be obtained from the source function by transformation of the frequency $\omega$.

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## 2. Uniqueness

Proposition 2.1. If $S$ is a source function for (1), then the integral
(2) $E(y, x, t) \doteq \frac{1}{2 \pi} \int_{\Gamma} \exp (-\imath \omega t) S(y, x ; \omega) d \omega ; \quad \Gamma \doteq\left\{\omega: \operatorname{Im} \omega=\omega_{0}+\tau\right.$ and $\left.\tau>0\right\}$,
is the forward fundamental solution of the operator $\square \doteq \Delta_{x}-\mathbf{n}^{2} \partial_{t}^{2}+\mathbf{m}$.
Proof. The integral (2) converges in the space of tempered distributions to a fundamental solution of the wave operator. On the other hand, it does not change if we replace $\Gamma$ by $\Gamma+\tau \imath, \tau>0$. This gives the factor $\exp (-\tau t)$ in an estimate of $E$ : which implies $E=0$ for $t<0$, i.e. $E$ is the forward fundamental solution of the wave operator.

Corollary 2.2. Any two source functions coincide.
Proof. This follows from the uniqueness of the forward fundamental solution. To check the uniqueness, we note that the difference of any two forward fundamental solutions vanishes for $t<0$ and consequently for all $t$ by compactness of the dependence domain.

## 3. The main result

We assume the following conditions:
(*) $\mathbf{n} \in C^{3}(X), \mathbf{n}>0$, and $\mathbf{n}=\mathbf{n}_{0}$ is constant in $X \backslash D$. where $D$ is a compact set in the plane $X$;
(**) the metric $g=\mathbf{n}^{2} d s^{2}$ is nontrapping in the plane, i.e. any geodesic curve $\gamma$ in $D$ quits the set $D$ in finite time:
$(* * *) \mathbf{m} \in C^{0}(X)$, and $\mathbf{m}=\mathbf{m}_{0}$ is constant in $X \backslash D$.
Denote by $t(\cdot, \cdot)$ the distance function in the space $(X, g)$ and by $t_{0}(y, x)$ the distance for the metric $g_{0}=\min \left\{\mathbf{n}^{2}(x) \cdot \mathbf{n}_{0}^{2}\right\} d s^{2}$.

Theorem 3.1. Under the conditions $(*),(* *)$ and $(* * *)$ there exist a positive $\omega_{0}$, and a source function $S(y, x ; \omega)$ in $X \times X \times \mathbf{C}_{+}\left(\omega_{0}\right)$ for (1) that satisfies the inequality

$$
\begin{equation*}
\int_{t_{0}(y, x) \leq \tau}\left|\exp \left(\operatorname{Im} \omega t_{0}(y, x)\right) S(y, x: \omega)\right|^{2} d x \leq \frac{1}{|\omega|}\left(C_{1} \tau+C_{2} \tau^{5}\right) \tag{3}
\end{equation*}
$$

uniformly for $y \in X$ and $\tau \geq 0$, that is holomorphic and continuous with respect to $\omega$ in the closure of $\mathbf{C}_{+}\left(\omega_{0}\right)$ for $x \neq y$ and in the sense of $L_{2 . l o c}(X)$ near the diagonal $x=y$.

Note that a result of N. Burq [4] contains the estimates $\|S\|=O\left(\omega^{-1 / 2}\right)$ for local operator $L_{2}$-norms in spaces $X$ of arbitrary dimension for $\omega \rightarrow \infty$ near the real axis.

Given a forward fundamental solution $E$ for the wave operator in a medium with velocity $\mathbf{n}^{-1}$, we could find the source function by means of the inverse Fourier transform of $E$, provided we can control the growth of $E$ as $t \rightarrow \infty$. A local construction of the fundamental solution $E$ for second order operators with analytic coefficients was done by Hadamard [7]. The case of smooth coefficients was considered by Sobolev [9]. Duistermaat and Hörmander [5] have constructed the forward propagator for arbitrary strictly hyperbolic operators by the method of Fourier integral operators. However, this approach is not easy to implement directly, even in the situation of Theorem 3.1, since there is not enough information on the growth of the fundamental solution as $t \rightarrow \infty$. We here apply a more involved method. First, we construct a forward parametrix $P$ for the wave equation. For this we choose a consistent chain of singular functions starting from Hadamard's fundamental solution. Our construction is parallel to that of [5]. but we use simple singular functions rather than Fourier integrals. The time Fourier transform of $P$ will be a parametrix for (1). Finding the correction term, we complete the construction.

Remark. The condition $\operatorname{dim} X=2$ helps to keep the volume of this paper limited. No essential difficulty appears in the general case, except for that a more detailed analysis of the energy in the caustic area is necessary.

Corollary 3.2. The forward fundamental solution for arbitrary $\tau$ satisfies the inequality

$$
\int_{t_{0}(y, x) \leq \tau} \int_{\mathbf{R}}\left|E_{\varepsilon}(y, x, t)\right|^{2} d t d x \leq C_{\phi}\left(C_{1} \tau+C_{2} \tau^{5}\right) \log \varepsilon
$$

where

$$
E_{\varepsilon}(y, x, t) \doteq \frac{1}{\varepsilon} \int_{\mathbf{R}} \exp \left(-\omega_{0} s\right) E(y, x, s)_{\varphi}\left(\frac{t-s}{\S}\right) d s
$$

and $\varphi$ is an arbitrary function in $\mathbf{R}$ with compact support such that

$$
\int_{0}^{\infty}|\widehat{\rho}( \pm \exp \sigma)|^{2} d \sigma<\infty
$$

Proof. By Parseval's theorem,

$$
\int_{\mathbf{R}}\left|E_{\S}\right|^{2} d t=\frac{1}{2 \pi} \int_{\mathbf{R}}\left|\hat{\varphi}\left(\varepsilon u^{\prime}\right) S\left(y \cdot x: \omega+\imath \omega_{0}\right)\right|^{2} d \omega
$$

and we apply (3).
Corollary 3.3. The source function has a symmetric kernel (the reciprocity property), i.e. $S(y, x ; \omega)=S(x . y ; \omega)$.

Proof. Take an arbitrary $\omega$, such that $\operatorname{Im} u>0$ : then for arbitrary $y$ and $x$

$$
\begin{aligned}
S(y, x ; \omega) & =\int_{X} S(y, z ; \omega) \delta(z-x) d z=\int_{X} S(y, z: \omega) \square_{u} S(x, z ; \omega) d z \\
& =\int_{X} S(x, z ; \omega) \square_{u} S(y, z ; \omega) d z=\int_{X} S(x, z ; \omega) \delta(z-y) d z=S(x \cdot y ; \omega)
\end{aligned}
$$

where $\square_{\omega} \doteq \Delta_{z}+\omega^{2} \mathbf{n}^{2}+\mathbf{m}$ and the function $S(y, z: \omega)$ is fast decreasing with respect to $z$ according to (3). The equation $S(y . x: \omega)=S(x . y ; \omega)$ is valid for real $\omega$ because of the uniqueness of the analytic continuation.

Corollary 3.4. The function $S(y, x ; w)$ is real for $v>w_{0}$.
Proof. The function $\widetilde{S}(y, x ; \omega) \doteq \bar{S}(y, x:-\bar{\omega})$ is also a source function for the real operator (1). By uniqueness it coincides with $S$.

Proposition 3.5. The source function satisfies the Sommerfeld radiation conditions:

$$
S(y, x ; \omega)=O\left(r^{-1 / 2}\right) \text { and }\left(\frac{\partial}{\partial r}-\imath \dot{\mathbf{n}_{0}}\right) S\left(y, x: \dot{)}=O\left(r^{-3 / 2}\right), \text { as } r \doteq|x-y| \rightarrow \infty\right.
$$

Proof. Apply the Helmholtz operator $\square_{\omega 0} \doteq \Delta_{x}+\omega^{2} \mathbf{n}_{0}^{2}+\mathbf{m}_{0}$ for the homogeneous medium

$$
\square_{\omega 0} S(y \cdot x: \omega)=\delta_{y}(x)+T(y \cdot x: \omega)
$$

where $T(y, x ; \omega) \doteq\left[\mathbf{m}_{0}-\mathbf{m}(x)+\omega^{2}\left(\mathbf{n}_{0}^{2}-\mathbf{n}^{2}(x)\right)\right] S(y \cdot x: \omega)$ and set

$$
\widetilde{S}(y, x ; w)=S_{0}\left(x-y ; u^{\prime}\right)+\int_{D} S_{0}(x-z: w) T(y . z: \dot{w}) d z
$$

where $S_{0}(x ; \omega)$ is the source function in the homogeneous medium. We have

$$
\square_{\omega 0} \widetilde{S}(y, x ; \omega)=\delta_{y}+T\left(y, x: \omega^{\prime}\right)
$$

It follows that $\square_{\omega 0} U(y, x ; \omega)=0$, where $U=S-\widetilde{S}$. Note that the function $\widetilde{S}$ has analytic continuation to the halfplane $\mathbf{C}_{+}(0)$ since the function $S_{0}$ has such a continuation. Apply Proposition 16.1 for the coefficient $\mathbf{p}_{0}=\sqrt{\omega^{2} \mathbf{n}_{0}^{2}+\mathbf{m}_{0}}$, where $|\omega|>\left|\mathbf{m}_{0}\right|^{1 / 2} \mathbf{n}_{0}^{-1}$, and $\operatorname{Im} \mathbf{p} \geq 0$. This gives $\widetilde{S}=O\left(|\omega|^{3 / 2} \exp \left(-\operatorname{Im} \mathbf{p}_{0}|x|\right)\right)$ for large $\omega^{\prime} ;$ this, together with (3), implies

$$
\begin{equation*}
\int_{t_{0}(y, x) \leq \tau}\left|\exp \left(\operatorname{Im} \omega t_{0}(y, x)\right) U(y, x ; \omega)\right|^{2} d x \leq|\dot{\nu}|^{3}\left(C_{1} \tau+C_{2} \tau^{5}\right) \tag{4}
\end{equation*}
$$

The Fourier-Laplace transform

$$
V(y, x, t) \doteq \int_{\Gamma} \exp (-\imath \omega t) U(y . x: \omega) d \omega
$$

is well defined for the curve $\Gamma \doteq\{\omega: \operatorname{Im} \omega=\tau\}$. By (4) the integral converges for arbitrary $\tau>0$ in the distribution sense and satisfies $\square_{0} V=0$. It decays to zero for $t<0$, as $\tau \rightarrow \infty$, but, on the other hand, does not depend on $\tau$. Therefore $V=0$ for $t<0$; hence $V=0$, by the uniqueness theorem for the wave operator. It follows that $S=\widetilde{S}$. We check that $\widetilde{S}$ satisfies the radiation condition. It is true for $S_{0}$ and also for the term $\int_{D} S_{0}(x-z ; \omega) T(y, z ; \omega) d z$, since the integral is taken over the compact set $D$.

Remark. For Schrödinger-type equations, the limiting absorption principle and the radiation conditions were studied by D. Eidus [6].

## 4. The diffusion equation

The diffusion equation in the optical tomography is

$$
\begin{equation*}
\frac{1}{\mathrm{c}} \frac{\partial \Phi}{\partial t}-\langle\nabla . \not \approx \nabla\rangle \Phi+\mu \Phi=q \tag{5}
\end{equation*}
$$

where $\Phi=\Phi(x, t)$ - the photon density, $q$ the density of the source. $\mu=\mu(x)$ is the absorption coefficient measured in the unit metre ${ }^{-1}, x(x)>0$ is the inverse diffusion coefficient measured in metres, and $\mathbf{c}$ is the light velocity measured in metre/second. This is the $P_{1}$-approximation to the transport equation, [1]. In frequency domain it appears as

$$
-\langle\nabla \cdot x \nabla\rangle \hat{\Phi}+\mu \widehat{\Phi}+\imath \frac{\lambda}{\mathbf{c}} \widehat{\Phi}=\hat{q} .
$$

For the unknown function $\Psi=\varkappa^{1 / 2} \Phi$ the equation

$$
\Delta \Psi+\left(-\imath \frac{\lambda}{\mathbf{c} \varkappa}+\mathbf{m}\right) \Psi=-\varkappa^{-1 / 2} q
$$

holds, where $\mathbf{m} \doteq-\varkappa^{-1 / 2} \Delta \varkappa^{1 / 2}-\varkappa^{-1} \mu$. It takes the form (1) for $\mathbf{n}=(\mathbf{c} \varkappa)^{-1 / 2}$, $\omega^{2}=-\imath \lambda$, where $\lambda$ is endowed by dimension second ${ }^{-2}$. Define

$$
S_{D}(y, x ; \lambda) \doteq-\varkappa^{-1 / 2} S_{H}(y, x: \sqrt{-\imath \lambda})
$$

where $S_{H}$ is the source function of $(1)$ and $\operatorname{Im} \omega \doteq \operatorname{Im} \sqrt{-\imath \lambda}>0$.
Corollary 4.1. If the coefficients $\mu \in C^{0}(X)$ and $\varkappa \in C^{3}(X)$ are constant in $X \backslash D$ for a compact set $D$ and $\mathbf{n}=(\mathbf{c} \varkappa)^{-1 / 2}$ satisfies the condition $(* *)$, then there exists a function $S_{D}(y, x ; \lambda)$ defined in $X \times X \times\left(\mathbf{C} \backslash \mathbf{R}_{+}\right)$that possesses the following properties:
(i) it satisfies

$$
\left(-\langle\nabla, \varkappa \nabla\rangle+\mu+\imath \frac{\lambda}{\mathbf{c}}\right) S_{D}(y, \cdot ; \lambda)=\delta(\cdot-y)
$$

(ii) it is holomorphic in $\lambda$, and
(iii) it decreases fast as $t_{0}(y, x) \rightarrow \infty$ and the inequality

$$
\int_{t_{0}(y, x) \leq \tau}\left|\exp \left(\sqrt{|\lambda|-\operatorname{Im} \lambda} t_{0}(y, x)\right) S_{D}(y \cdot x ; \lambda)\right|^{2} d x \leq \frac{1}{|\lambda|^{1 / 2}}\left(C_{1} \tau+C_{2} \tau^{5}\right)
$$

holds for $\tau>0$, where $t_{0}$ is the distance function as in Theorem 3.1.
Corollary 4.2. The function

$$
E(y, x, t)=\frac{1}{2 \pi} \int_{\operatorname{Im} \lambda=\tau<0} S_{D}(y \cdot x ; \lambda) \exp (\imath \lambda t) d \lambda
$$

is a fundamental solution of the diffusion equation (5).

## 5. The parametrix of the wave equation

First, we construct a forward parametrix $P=P(y, x, t)$ for the wave operator $\square \doteq \Delta_{x}-\mathbf{n}^{2} \partial_{t}^{2}+\mathbf{m}$ in space-time $X \times \mathbf{R}$. We shall use the coordinates $x, t ; \xi, \tau$ in the phase space $\Phi \doteq(X \times \mathbf{R}) \times(X \times \mathbf{R})^{*}$, where $\xi$ and $\tau$ are conjugate to the space-time coordinates $x$ and $t$. This means that $\alpha \doteq \xi d x+\tau d t$ is the canonical contact form
in $\Phi$, where $\xi d x=\xi_{1} d x_{1}+\xi_{2} d x_{2}$. The function $\mathbf{n}^{2}(x) \tau^{2}-|\xi|^{2}$ is the symbol of this operator. Consider the Hamiltonian flow generated by $h(x: \xi, \tau) \doteq \frac{1}{2}\left(\mathbf{n}^{2}(x) \tau^{2}-|\xi|^{2}\right)$ :

$$
\begin{equation*}
\frac{d x}{d r}=-\xi, \quad \frac{d t}{d r}=\mathbf{n}^{2}(x), \quad \frac{d \xi}{d r}=-\mathbf{n}(x) \nabla \mathbf{n}(x) . \quad r \geq 0 ; \tau=1 \tag{6}
\end{equation*}
$$

Fix a point $y \in X$ and take all the solutions with initial data

$$
x(0)=y, \quad t(0)=0, \quad \xi(0) \in T_{y}^{*}(X) \quad \text { and } \quad|\xi(0)|=\mathbf{n}(y) .
$$

Let $\Lambda_{y}$ be the union of all trajectories of the flow. This is a smooth surface in $\Phi$. We have $|\xi|=\mathbf{n}(x)$ and the contact form $\alpha$ vanishes in $\Lambda_{y}$, by Jacobi's theorem. This means that the conic set in $\Phi$ generated by $\Lambda_{y}$ is a Lagrange manifold.

Let $\pi: \Phi \rightarrow X$ be the natural projection; the image of a solution of (6) is called a ray. It is a geodesic for the metric $g$. The set $L_{y} \doteq \pi\left(\Lambda_{y}\right)$ is the union of rays starting at $y$; it is called the front of the wave which starts at $y$. The set $K_{y} \doteq$ $L_{y}+\{(0, t): t \geq 0\}$ is called the future conoid of $y$ : the boundary $\partial K_{y}$ is the first wave front. Consider the union $K \doteq \bigcup_{y \in X}\left(\{y\} \times K_{y}\right)=\{(y, x, t): t \geq t(y, x)\}$.

Lemma 5.1. There exists a kernel $P$ in $X \times X \times \mathbf{R}_{+}$supported by $K$ such that its Fourier transform in $\mathbf{C}_{+}(0)$ satisfies the inequalities

$$
\begin{align*}
& \int_{t(y, x) \leq \tau} \exp (2 \operatorname{Im} \omega t(y, x))|\hat{P}(y, x ; \omega)|^{2} d x \leq \frac{1}{\left|\omega^{\prime}\right|}\left(C_{1} \tau+C_{2} \tau^{5}\right)  \tag{7}\\
& \int_{t(y, x) \leq \tau} \exp (2 \operatorname{Im} \omega t(y, x))|\widehat{Q}(y, x ; \omega)|^{2} d x \leq \frac{1}{\left|\omega^{\prime}\right|^{3}}\left(C_{1} \tau+C_{2} \tau^{5}\right) \tag{8}
\end{align*}
$$

where $Q \doteq \square P-\delta_{y, 0}$.
The notation $\hat{A}$ is used for the time Fourier transform:

$$
\hat{A}(\omega) \doteq F(A)=\frac{1}{2 \pi} \int_{\mathbf{R}} \exp (-\imath \omega t) A(t) d t
$$

Proof. We construct the kernel $P$ in several steps. We will let $e$ be a real smooth function with compact support in $\mathbf{R}$ that is equal to 1 in a neighbourhood of the origin; it need not be the same in all steps.

Step 1. Choose $\tau>0$ so small that the geodesic coordinates are defined in the neighbourhood $U_{0} \doteq\left\{(y, x): t(y, x) \leq \tau_{0}\right\}$ of the diagonal in $X \times X$. We define the parametrix in $U_{0}$ by Hadamard's method:

$$
\begin{align*}
A_{0}(y, x, t) & \doteq \theta(t-t(y, x)) e(t-t(y, x))\left(t^{2}-t^{2}(y, x)\right)^{-1 / 2} a(y, x, t)  \tag{9}\\
a(y, x, t) & =a_{0}(y, x)+\left(t^{2}-t^{2}(y, x)\right) a_{1}(y, x)
\end{align*}
$$

Here $t(y, x)$ denotes the distance from $y$ to $x$ in the metric space $(X, g)$; it is a smooth function in $U_{0}$. The amplitude $a(y, x)$ is a smooth function with respect to the polar coordinates centered in $y$ such that

$$
\begin{equation*}
\square A_{0}(y, x, t)=\delta(x-y)+\theta(t-t(y \cdot x))\left(t^{2}-t^{2}(y \cdot x)\right)^{1 / 2} b(y, x, t)+B(y, x, t) \tag{10}
\end{equation*}
$$

for some smooth functions $b$ and $B$ : where $\theta(t)=1$ for $t \geq 0$ and $\theta(t)=0$ otherwise.
Step 2. We extend Hadamard's construction to $D_{0} \times D_{0}$, where $D_{0}$ is a compact neighbourhood of $D$, by choosing a consistent chain of singular functions $A_{0}$ starting with $A_{0}$. The union $\Lambda \doteq \bigcup_{X}\left(\{y\} \times \Lambda_{y}\right)$ is a smooth manifold in $X \times \Phi$. Consider the mapping $\pi: \Lambda \rightarrow X \times \Phi \rightarrow X \times X$, that equals the composition of projections. For an arbitrary $y \in X$ we denote by $\pi_{y}: \Lambda_{y} \rightarrow X$ the restriction of $\pi$. The rank of $\pi_{y}$ is equal to 1 or 2 in each point. since the projection $\pi_{y}: \gamma \rightarrow X$ is an immersion for an arbitrary trajectory $\gamma$ of the flow (6).

Take a point $(y, \lambda) \in \Lambda$ such that rank $d \pi_{y}(\lambda)=2$ and choose a generating function for the germ of $\Lambda$ in the form $\phi(y, x, t)=t-\hat{\varphi}(y \cdot x)$. where $\hat{\rho}$ is an eikonal function. We take the singular function

$$
\begin{equation*}
A^{(2)}(y, x, t) \doteq \phi(y, x, t)_{ \pm}^{-1 / 2} a(y, x, t) \tag{11}
\end{equation*}
$$

that depends on the phase $\sigma$ and the amplitude $a$. Here and later we use the notation $s_{ \pm}^{\lambda} \doteq|s|^{\lambda}$ for $\pm s \geq 0$ and $s_{ \pm}^{\lambda}=0$ otherwise. The amplitude has the form $a(y, x, t)=$ $e(\phi(y, x, t))\left[a_{0}(y, x)+2 \phi(y . x, t) a_{1}(y, x)\right]$. where $a_{0}$ and $a_{1}$ are some smooth functions and $e(\tau) \in \mathcal{D}(\mathbf{R})$ is a function that is equal to 1 in a neighbourhood of the origin. They are subjected to the system of transport equations

$$
\begin{gather*}
2\left\langle\nabla \varphi, \nabla a_{0}\right\rangle+\Delta(\varphi) a_{0}=0, \\
2\left\langle\nabla \varphi, \nabla a_{1}\right\rangle+[\Delta(\varphi)+\mathbf{m}] a_{1}=\Delta a_{0} . \tag{12}
\end{gather*}
$$

Note that the function $A_{0}$ (see (9)) is of type $A^{(2)}$ everywhere in the set $\{(y, x, t)$ : $x \neq y\}$. We take the kernel $\phi(y, x, t)_{+}^{-1 / 2}$ in (11). if the point $(x, t=\varphi(y, x))$ belongs to the first front, i.e. to the boundary of $K_{y}$ and also to any regular point of $L_{y}$ such that the ray $\gamma(y, x, t)$ from ( $y .0$ ) to ( $x . t$ ) has even Morse index. The kernel $\phi(y, x, t)_{-}^{-1 / 2}$ is used for any regular point of $L_{y}$ with odd Morse index of the corresponding ray $\gamma(y, x, t)$. We have

$$
\begin{equation*}
\square A^{(2)}=o_{ \pm}^{1 / 2} e(o) \Delta a_{1}+B \tag{13}
\end{equation*}
$$

where $B$ is a smooth function.

Lemma 5.2. The function $e$ can be chosen in such a way that $\operatorname{supp} A^{(2)} \subset K$.
Proof. The choice of $e$ is not important in the case of the kernel $\delta_{+}^{-1 / 2}$, since $\operatorname{supp} \phi_{+}^{-1 / 2} \subset K$ for an arbitrary eikonal function $\gamma$. The kernel $\sigma_{-}^{-1 / 2}$ appears only if the Morse index of the corresponding geodesic $\gamma_{j}(y . x . t)$ is odd. By Jacobi's theorem (see e.g. [3]) this implies that this geodesic is not the shortest path from $y$ to $x$. Therefore the point $(y, x, \varphi(y, x)) \in L_{y}$ does not belong to the boundary of $K_{y}$. We can choose the function $e$ in such a way that the function $e(\phi(y, x, t))$ vanishes on $\partial K$.

Step 3. Take an arbitrary point $(y, \lambda) \in \Lambda$ such that rank $d \pi_{y}(\lambda)=1$. We have $d t \neq 0$ in $\Lambda$ and there exists a Euclidean coordinate system $x=(u, v) . t$ in a neighbourhood of $\pi(\lambda)$ such that $\xi d x \doteq \eta d u+\theta d v$ for the dual coordinates $\xi=(\eta, \theta)$ and $(y, u ; \theta)$ is a local coordinate system in $\Lambda$. We have $t=t(u . \theta)$ and $v=v(u, \theta)$ in $\Lambda$ for some smooth functions. The phase function

$$
\phi(y, x, t ; \theta)=t-\varphi(y \cdot x: \theta), \quad \varphi(y \cdot x: \theta)=t(u . \theta)+\theta(v(u \cdot \theta)-v)
$$

generates a neighbourhood $\Lambda^{\prime}$ of $\lambda$ in $\Lambda$, i.e.

$$
\Lambda^{\prime}=\left\{(y, x, t ; \xi, \tau): \eta=\phi_{u}^{\prime}(y, x, t ; \theta), \quad \phi(y, x, t ; \theta)=0 \text { and } o_{\theta}^{\prime}(y, x, t ; \theta)=0\right\}
$$

is a neighbourhood of $\lambda$ in $\Lambda$. Indeed, the equations $t_{\theta}^{\prime}+\theta v_{\theta}^{\prime}=0, t_{u}^{\prime}+\theta v_{u}^{\prime}+\eta=0$ hold in $\Lambda$, since $\alpha=0$.

We shall use the notation $p_{2}(s)=s \log (s+0 \imath)-s, p_{1}(s)=\log (s+0 \imath)$ and $p_{k-1}=p_{k}^{\prime}$ for $k=-1,0,1,2$. Take a singular function of the form

$$
\begin{equation*}
A^{(1)}(y, x, t)=\int_{\mathbf{R}} e(\phi) \operatorname{Re}\left[p_{0}(\phi) a_{0}+p_{1}(\phi) a_{1}\right] d \theta \tag{14}
\end{equation*}
$$

where the phase function $\phi=t-\varphi$ satisfies the generalized eikonal equation $\mathbf{n}^{2}-$ $|\nabla \varphi|^{2}=\varphi_{\theta}^{\prime} \psi$, where $\psi=\psi(x, \theta)$ is a smooth function. It exists and is unique, since the left-hand side vanishes as $\varphi_{\theta}^{\prime}=v(u, \theta)-v=0$. The amplitude functions $a_{k}=$ $a_{k}(y, x ; \theta), k=0,1$, are smooth and have proper supports (that is the projections supp $a_{k} \rightarrow X \times X,(y, x ; \theta) \mapsto(y, x)$ are proper $)$. They are subjected to the transport equations

$$
\begin{align*}
2\left\langle\nabla \varphi, \nabla a_{0}\right\rangle+\left(\psi a_{0}\right)_{\theta}^{\prime}+\Delta(\varphi) a_{0} & =\varphi_{\theta}^{\prime} b_{0} \\
2\left\langle\nabla \varphi, \nabla a_{1}\right\rangle+\left(\psi a_{1}\right)_{\theta}^{\prime}+[\Delta(\varphi)+\mathbf{m}] a_{1} & =2 \Delta a_{0}+\varphi_{\theta}^{\prime} b_{1}-b_{0 \theta}^{\prime} \tag{15}
\end{align*}
$$

for some smooth functions $b_{0}$ and $b_{1}$. Note that (12) and (15) are the only points where the coefficient $\mathbf{m}$ contributes to our construction.

Lemma 5.3. If the functions $a_{0}$ and $a_{1}$ are supported by $X \times X \times \Theta$, where $\Theta$ is a sufficiently small neighbourhood of $\theta_{0}=\theta(\lambda)$, we can choose a function e in such a way that $\operatorname{supp} A^{(1)} \subset K$.

Proof. We argue as in Lemma 5.2. The Morse index of the ray $\gamma(y, x, t)$ changes just when the curve passes a singular point $\pi(\lambda) \in L_{y}$ such that rank $d \pi_{y}(\lambda)=1$. This point does not belong to $\partial K_{y}$, due to Jacobi's theorem. Therefore there is a space for choosing an appropriate function $e$.

Lemma 5.4. We have

$$
\begin{equation*}
\square A^{(1)}(x, t)=\int_{\mathbf{R}} e(\phi) \operatorname{Re}\left(p_{1}(\phi)\left(\Delta a_{1}(x, \theta)-b_{1 \theta}^{\prime}(x, \theta)\right)\right) d \theta+B(x, t), \tag{16}
\end{equation*}
$$

where $B$ is a smooth function.
Proof. We calculate

$$
\begin{aligned}
\square A^{(1)}(x, t)= & \int_{\mathbf{R}} \operatorname{Re} p_{-2}(\phi)\left(\left(\nabla_{\varphi}(x, \theta)\right)^{2}-\mathbf{n}^{2}(x)\right) a_{0} d \theta \\
& +\int_{\mathbf{R}} \operatorname{Re} p_{-1}(\phi)\left[\left(\left(\nabla_{\varphi}\right)^{2}-\mathbf{n}^{2}\right) a_{1}-2\left\langle\nabla_{\varphi}, \nabla a_{0}\right\rangle-\Delta \varphi a_{0}\right] d \theta \\
& +\int_{\mathbf{R}} \operatorname{Re} p_{0}(\phi)\left[-2\left\langle\nabla \varphi, \nabla a_{1}\right\rangle-\Delta \varphi a_{1}+(\Delta+\mathbf{m}) a_{0}\right] d \theta \\
& +\int_{\mathbf{R}} \operatorname{Re} p_{1}(\phi) \Delta a_{1} d \theta+B .
\end{aligned}
$$

The term $B$ contains derivatives of $e(\phi)$ and is a smooth function. We integrate the first term by parts using the eikonal equation.

$$
\int_{\mathbf{R}} p_{-2}(\phi)\left((\nabla \varphi)^{2}-\mathbf{n}^{2}\right) a_{0} d \theta=-\int_{\mathbf{R}} p_{-2}(\phi) \varphi_{\theta}^{\prime} v a_{0} d \theta=-\int_{\mathbf{R}} p_{-1}(\phi)\left(\psi a_{0}\right)_{\theta}^{\prime} d \theta
$$

Combining this with the second term we obtain

$$
\begin{aligned}
\int_{\mathbf{R}} p_{-1}(\phi)\left(\left((\nabla \varphi)^{2}-\mathbf{n}^{2}\right) a_{\mathbf{1}}\right. & \left.-2\left(\nabla \varphi, \nabla a_{0}\right\rangle-\Delta \varphi a_{0}-\left(v a_{0}\right)_{\theta}^{\prime}\right) d \theta \\
& =-\int_{\mathbf{R}} p_{-1}(\phi) \varphi_{\theta}^{\prime}\left(\psi a_{1}+b_{0}\right) d \theta=-\int_{\mathbf{R}} p_{0}(\phi)\left(\psi a_{1}+b_{0}\right)_{\theta}^{\prime} d \theta
\end{aligned}
$$

This together with the third term gives

$$
\begin{aligned}
\int_{\mathbf{R}} p_{0}(\phi)\left[-2\left\langle\nabla \varphi, \nabla a_{1}\right\rangle-\Delta \varphi a_{1}+(\Delta+\mathbf{m}) a_{0}-\left(v a_{1}+b_{0}\right)_{\theta}^{\prime}\right] d \theta & =-\int_{\mathbf{R}} p_{0}(\phi) \varphi_{\theta}^{\prime} b_{1} d \theta \\
& =-\int_{\mathbf{R}} p_{1}(\phi) b_{1 \theta}^{\prime} d \theta
\end{aligned}
$$

and so on.

Step 4. Take a bounded convex set $D_{0}$ that contains a neighbourhood of $D$; we shall specify it later.

Lemma 5.5. There exist a finite system of open sets $\left\{U_{i} \subset X \times X: 0 \leq i \leq i_{*}\right\}$ that covers $D \times D_{0}$, a system of open sets $\left\{\Lambda_{\alpha}: 0 \leq \alpha \leq \alpha_{*}\right\}$ in $\Lambda$, a mapping $i=i(\alpha)$ for $0 \leq \alpha \leq \alpha_{*}$ such that $\pi\left(\Lambda_{\alpha}\right) \subset U_{i(\alpha)}$ for any $\alpha \geq 0$, and for any $\alpha$ a function $A_{\alpha}$ in $U_{i(\alpha)}$ such that
(i) $U_{0}$ is a neighbourhood of the diagonal in $D \times D_{0}$ and $A_{0}$ is defined by (9):
(ii) for any $\alpha>0, A_{\alpha}$ is a function of form (11) or (14), where the phase function $\phi=\phi_{\alpha}(y, x, t)$ and $\phi=\phi_{\alpha}(y, x, t ; \theta)$. respectively, generates $\Lambda_{\alpha}$;
(iii) for arbitrary $j$ and $k$ the function

$$
B_{j k} \doteq \sum_{i(\alpha)=j} A_{\alpha}-\sum_{i(\beta)=k} A_{\beta}
$$

can be represented in $U_{j} \cap U_{k}$ as

$$
\begin{align*}
B_{j k}(y, x, t)= & \sum\left(\phi_{\gamma}(y, x, t)\right)_{ \pm}^{1 / 2} b_{\gamma}(y, x, t) \\
& +\sum \int_{\mathbf{R}} \operatorname{Re}\left(p_{2}\left(\phi_{\gamma}(y, x, t ; \theta)\right) b_{\gamma}(y, x . t: \theta)\right) d \theta \tag{17}
\end{align*}
$$

where the sums are taken over all $\gamma$ such that $i\left(\gamma_{i}\right)=j, k$; and $b_{\gamma}$ are smooth functions with proper supports.

Proof. The function (9) is well defined in $U_{0}$. The projection $\pi: \Lambda_{0} \rightarrow U_{0}$ is bijective for $\Lambda_{0}=\pi^{-1}\left(U_{0}\right) \cap \Lambda$. Next, we construct the functions $A_{\alpha}$ for all $\alpha=$ $1, \ldots, \alpha_{1}$ such that $\Lambda_{\alpha}$ has a nonempty intersection with $\Lambda_{0}$. Then we find the functions $A_{\beta}$ for all $\beta=\alpha_{1}+1, \ldots, \alpha_{2}$ such that $\Lambda_{3} \cap \Lambda_{\alpha} \neq \emptyset$ at least for one $\alpha \leq \alpha_{1}$ and so on. The initial conditions for the amplitudes are defined by means of (iii). Consider one step of the continuation in detail. If the functions $A_{\alpha}$ and $A_{\beta}$ are of type $A^{(2)}$, and the phase functions $\varphi_{\alpha}$ and $\psi_{3}$ coincide in $\pi\left(\Lambda_{\alpha}\right) \cap \pi\left(\Lambda_{3}\right)$. We can take a smooth continuation of amplitudes $a_{\alpha k}$ to $a_{3 k}$ for $k=0,1$ preserving the equation (12).

Assume that $A_{\alpha}$ and $A_{\beta}$ are of type $A^{(1)}$. Suppose that there exists a point $(y, \lambda) \in \Lambda_{\alpha} \cap \Lambda_{\beta}$ such that rank $d \pi_{y}(\lambda)=1$. Then there exists a coordinate system $(v, \eta)$ in a neighbourhood of $\lambda$ such that $\theta=\theta(\eta) . x=x(v, \eta)$,

$$
\varphi_{\alpha}(x(v, \eta) ; \theta(\eta))= \pm \varphi_{\beta}(v ; \eta) \quad \text { and } \quad \frac{a_{\alpha}(x(v, \eta) ; \theta(\eta))}{\left|\eta_{\theta}^{\prime}\right|}=a_{3}(v ; \eta)
$$

The integral $A_{\alpha}$ is transformed to $A_{\beta}$ by the coordinate change. This transformation also does not contribute to the sum $B_{U, V}$. If there is no such point $\lambda$, we can transform $A_{\alpha}$ to $A_{\beta}$ through a function of type $A^{(2)}$.

The most complicated case is the transformation of a function $A_{\alpha}$ of type $A^{(1)}$ to a sum of functions $A_{\beta}$ of type $A^{(2)}$. We assume for simplicity that $a_{1}=0$. Take an arbitrary point $(y, \lambda) \in \Lambda_{\alpha} \cap \Lambda_{\beta}$ such that $\operatorname{rank} d \pi_{y}(\lambda)=2$ and set $x=\pi_{y}(\lambda)$. We have $\theta=\theta_{\beta}(x)$ in a neighbourhood of $x$, due to the rank condition and $\partial_{\theta}^{2} \varphi_{\alpha}\left(x ; \theta_{3}(x)\right) \neq 0$ since $\phi_{\alpha}$ is a nondegenerate phase function. Therefore we have

$$
\begin{equation*}
\varphi_{\alpha}(x ; \theta)=\varphi_{3}(x) \pm \eta^{2} . \quad \text { where } \hat{\gamma}_{3}(x)=\hat{\varphi}_{\alpha}\left(x ; \theta_{3}(x)\right) \tag{18}
\end{equation*}
$$

for a smooth function $\eta=\eta(x, \theta)$ defined in a neighbourhood of $\left(y, \theta_{3}(x)\right)$ such that $\eta_{\theta}^{\prime} \neq 0$, where the phase function $o_{3}(x . t)=t-\gamma_{3}(x)$ generates $\Lambda_{3}$. By means of a partition of unity we write the above integral as a sum of integrals of the form

$$
B_{\beta}(x, t)=\int_{\mathbf{R}} \operatorname{Re} p_{0}\left(o_{3} \pm \eta^{2}\right) b(x . t ; \eta) d \eta
$$

where $b=a\left|\theta_{\eta}^{\prime}\right|$ is a smooth function with proper support. Write $b(x, t ; \eta)=b(x, t ; 0)+$ $c(x, t ; \eta) \eta$ for a smooth function $c$ and get

$$
\begin{align*}
B_{\beta}(x, t) & =\operatorname{Re}\left[b(x, t ; 0) \int_{\mathbf{R}} p_{0}\left(o_{3} \pm \eta^{2}\right) d \eta\right]+\int_{\mathbf{R}} \operatorname{Re} p_{0}\left(o_{3} \pm \eta^{2}\right) c(x, t ; \eta) \eta d \eta \\
9) & =\pi \operatorname{Re}\left[b(x, t ; 0)\left(\left(o_{3}\right)_{+}^{-1 / 2} \mp \imath\left(o_{3}\right)_{-}^{-1 / 2}\right)\right]-\frac{1}{2} \int_{\mathbf{R}} \operatorname{Re} p_{-1}\left(o_{\alpha}\right) c_{\eta}^{\prime}(x . t ; \eta) d \eta . \tag{19}
\end{align*}
$$

The first term is the sum of functions $A_{3}$ of type $A^{(2)}$. The second term is of type $A^{(1)}$; we denote it $B_{\alpha \beta}$. We have $A_{\alpha}^{(1)}=\sum\left(A_{3}^{(2)}+B_{\alpha 3}\right)$, where the sum is taken over all $\beta$ such that $\Lambda_{\alpha} \cap \Lambda_{\beta} \neq \emptyset$ and

$$
B_{\alpha \beta}(x, t)=\int_{\mathbf{R}} \operatorname{Re} p_{2}\left(o_{3}(x, t) \pm \eta^{2}\right) b_{3}(x, t ; \eta) d \eta
$$

for a smooth function $b_{\beta}$ with proper support.
We check that the system (15) for $\varphi_{0}$ is consistent with (12) for $\varphi_{3}$ in the hypersurface $H \doteq\left\{(x ; \eta):\left(\varphi_{\alpha}\right)_{\theta}^{\prime}(x ; \eta)=0\right\}$. By $(18),\left(\mathcal{F}_{\alpha}\right)_{\theta}^{\prime}= \pm 2 \eta \eta_{\theta}^{\prime}$, where $\eta_{\theta}^{\prime} \neq 0$, and we have $\nabla \varphi_{\alpha}=\nabla \varphi_{\beta}$ in $H=\{(x, \theta): \eta(x, \theta)=0\}$. From the generalized eikonal equation and (18) we also have

$$
\begin{aligned}
\left(\varphi_{\alpha}\right)_{\theta}^{\prime} \psi & =\mathbf{n}^{2}-\left|\nabla \varphi_{\alpha}\right|^{2}=\mathbf{n}^{2}-\left|\nabla \varphi_{3}\right|^{2} \mp 4 \eta\left\langle\nabla \varphi_{3}, \nabla \eta\right\rangle-4 \eta^{2}|\nabla \eta|^{2} \\
& \left.=-\left.4 \eta\left| \pm\left\langle\nabla \varphi_{\beta}, \nabla \eta\right\rangle+\eta\right| \nabla \eta\right|^{2}\right]
\end{aligned}
$$

which yields $\psi=-2\left(\eta_{\theta}^{\prime}\right)^{-1}\left\langle\nabla \varphi_{\beta}, \nabla \eta\right\rangle=-2\left(\eta_{\theta}^{\prime}\right)^{-1}\left\langle\nabla \varphi_{\alpha} . \nabla \eta\right\rangle$ in $H$. Therefore the field $2\left\langle\nabla \varphi_{\alpha}, \nabla\right\rangle+\psi \partial_{\theta}=2\left\langle\nabla \varphi_{\alpha}, \nabla\right\rangle-2\left(\eta_{\theta}^{\prime}\right)^{-1}\left\langle\nabla \varphi_{\alpha}, \nabla \eta\right\rangle \partial_{\theta}$ annihilates the function $\eta$ in $H$;
hence this field is defined in $H$. It coincides with $2\left\langle\nabla_{\nu_{3}}, \nabla\right\rangle$ on functions which do not depend on $\theta$, q.e.d. This implies that we can continue the amplitude functions from $\pi\left(\Lambda_{\alpha}\right)$ to $\pi\left(\Lambda_{\beta}\right)$ and vice versa. Finally we conclude that

$$
B_{j k}=\sum_{i(\alpha)=j} A_{\alpha}^{(1)}-\sum_{i(3)=k} A_{3}^{(2)}=\sum B_{\alpha^{3}}
$$

and Lemma 5.5 follows.
Step 5. Choose a family of functions $\left\{h_{i} \in \mathcal{D}\left(U_{i}\right)\right\}$ such that $\sum_{i} h_{i}(y, x)=1$ in $D \times D_{0}$. We set $P_{D} \doteq \sum_{\alpha} h_{i(\alpha)} A_{\alpha}$; the sum is finite and $A_{\alpha}$ are functions as in Lemma 5.5. We have

$$
\begin{equation*}
Q_{D} \doteq \square P_{D}=\sum_{j} \sum_{\alpha: i(\alpha)=j}\left[h_{j} \square A_{\alpha}+2\left\langle\nabla h_{j} . \nabla A_{\alpha}\right\rangle+\Delta h_{j} A_{\alpha}\right] . \tag{20}
\end{equation*}
$$

The term $\square A_{\alpha}$ is of the form (13) or (16). We show that the other terms cancel. up to a function of the form (17). Take an arbitrary point $(y, x) \in D \times D_{0}$ and fix an arbitrary index $j$. We have at this point

$$
\sum_{i(\alpha)=j}\left\langle\nabla h_{j}, \nabla A_{\alpha}\right\rangle=\sum_{i(3)=k}\left\langle\nabla h_{j} . \nabla A_{3}\right\rangle+\sum_{k}\left\langle\nabla h_{j} . \nabla B_{j k}\right\rangle .
$$

Take the sum with respect to $j$; the first sum on the right-hand side vanishes, since $\sum_{j} \nabla h_{j}=0$. The second sum is of the form (17) according to Lemma 5.5. Therefore the third sum in (20) is equal to the sum of $\Delta h_{j} B_{j k}$, which yields

$$
\begin{equation*}
\square P_{D}=\sum_{j} \sum_{i(\alpha)=j} h_{j} \square A_{\alpha}+2 \sum_{j . k}\left\langle\nabla h_{j} . \nabla B_{j k}\right\rangle+\sum_{j . k} \Delta h_{j} B_{j k} . \tag{21}
\end{equation*}
$$

## 6. The energy near the source point

Lemma 6.1. The following inequality holds for $\tau \leq \tau_{0}$,

$$
\int_{t_{0}(y, x) \leq \tau} \exp (2 \operatorname{Im} \omega t(y, x))\left(\left|\hat{A}_{0}(y, x: \omega)\right|^{2}+\left|\omega \square_{\omega} \hat{A}_{0}(y, x ; \omega)\right|^{2}\right) d x \leq C \frac{\tau}{|\omega|}
$$

where $\square_{\omega} \doteq \Delta+\omega^{2} \mathbf{n}^{2}+\mathbf{m}$.
Proof. We choose geodesic normal coordinates $z$ in the set $\left\{(y, x): t_{0}(y, x) \leq \tau_{0}\right\}$ such that $t^{2}(y, x)=|z|^{2} \doteq z_{1}^{2}+z_{2}^{2}$. We have by (9).

$$
\begin{align*}
\hat{A}_{0}(y, x ; \omega)= & a_{0}(y, z) \int_{t \geq|z|} \exp (\imath \omega t)\left(t^{2}-|z|^{2}\right)^{-1 / 2} e(t-|z|) d t  \tag{22}\\
& +a_{1}(y, z) \int_{t \geq|z|} \exp (\imath \omega t)\left(t^{2}-|z|^{2}\right)^{1 / 2} e(t-|z|) d t
\end{align*}
$$

By changing $s=t-|z|$, we can write $\hat{A}_{0}(y, x: \omega)=\exp (\imath w|z|)\left[a_{0}(y, z) I_{-}+a_{1}(y, z) I_{+}\right]$, where

$$
I_{ \pm} \doteq \int_{0}^{\infty} \exp (\imath \omega s)\left(s^{2}+2|z| s\right)^{ \pm 1 / 2} e(s) d s
$$

We estimate $I_{-}$for large $\omega$; the only singular point in the integrand is the origin of the ray $\{s \geq 0\}$. We assume that $e(\tau)=1$ for $|\tau| \leq 1$ and continue $e$ as a smooth function in $\mathbf{C}$, which equals 1 in the unit disc. Following the Laplace method, we replace the ray in (22) by the chain

$$
\begin{aligned}
\gamma & =\gamma_{1} \cup \gamma_{2} \cup \gamma_{3} \\
\gamma_{1} & =\left\{s=\tau \exp \left(\frac{1}{2} \pi-\arg \omega\right): 0 \leq \tau \leq 1\right\} \\
\gamma_{2} & =\left\{s=\exp (\imath \psi): 0 \leq|\psi| \leq \frac{1}{2} \pi-\arg \omega\right\} \\
\gamma_{3} & =\{s \geq 1\}
\end{aligned}
$$

and get $I_{-}=I_{1}+I_{2}+I_{3}$, where $I_{1}$ gives the main contribution:

$$
\begin{aligned}
\left|I_{1}(z)\right| \doteq\left|\int_{0}^{1} \exp (-\omega \tau)\left(-\tau^{2}+2 \imath|z| \tau\right)^{-1 / 2} d \tau\right| & \leq \int_{0}^{\infty} \exp (-\omega \tau)(2 \tau|z|)^{-1 / 2} d \tau \\
& =\left(\frac{\pi}{2|\omega z|}\right)^{1 / 2}
\end{aligned}
$$

We have $|s+2| z|s| \geq C(|z|+1)$ in $\gamma_{2} \cup \gamma_{3}$ and Re $\omega^{\prime} s \leq 0$. Integration by parts in $I_{2}+I_{3}$ yields $\left|I_{2}+I_{3}\right| \leq C|\omega|^{-1}(|z|+1)^{-1 / 2}$. For the integral $I_{+}$the same arguments yield the estimate $\leq C|z / \omega|^{1 / 2}$, which implies

$$
\int_{t_{0} \leq \tau}\left|\exp (\operatorname{Im} \omega t(y, x)) \hat{A}_{0}(y, x ; \omega)\right|^{2} d x \leq \frac{C}{|\dot{\omega}|} \int_{|z| \leq \tau} \frac{d x}{|z|} \leq C^{\prime} \frac{\tau}{|\omega|}
$$

We can estimate the kernel $\square_{w} \hat{A}_{0}$ in a similar way by means of (10).

## 7. The energy in the perturbed domain

We estimate the Fourier transforms of $P_{D}$ and of $Q_{D} \doteq \square P_{D}-\delta_{y, 0}$. For any $\lambda \in \mathbf{C}, \operatorname{Re} \lambda>0$ and arbitrary $b \in \mathcal{D}(\mathbf{R})$,

$$
\int_{\mathbf{R}} \exp (\imath \omega \tau) \tau_{ \pm}^{\lambda-1} b(\tau) d \tau=\Gamma(\lambda)(\imath \omega \mp 0)^{-\lambda}\left[b(0)+O\left(\frac{1}{\omega}\right)\right]
$$

consequently the Fourier transform of a function of type (11) equals

$$
\begin{aligned}
\hat{A}^{(2)}(x ; \omega) & =\int_{\mathbf{R}} \exp (\imath \omega t) A^{(2)}(x, t) d t=\int_{\mathbf{R}} \exp (\imath \omega t)(t-\varphi)_{ \pm}^{-1 / 2} e(t-\varphi) a(x) d t \\
& = \pm(\imath \pi \omega \mp 0)^{-1 / 2} \exp (\imath \omega \varphi(x))\left[a(x)+O\left(\frac{1}{\omega}\right)\right]
\end{aligned}
$$

Note that for the function $A_{\alpha}=A^{(2)}$ the phase $\varphi(x)$ coincides with the length of a geodesic from $(y, 0)$ to $(x, t)$. Therefore $\varphi(x) \geq t(y, x)$ and the exponential factor does not surpass $\exp (-\operatorname{Im} \omega t(y, x))$. By (13) the function $\square_{\omega} \hat{A}^{(2)}=F\left(\square A^{(2)}\right)$ has a similar structure; hence

$$
\begin{equation*}
\left|\hat{A}^{(2)}(x ; \omega)\right|+\left|\omega \square_{\omega} \hat{A}^{(2)}(x ; \omega)\right| \leq \frac{C}{|\omega|^{3 / 2}} \exp (-\operatorname{Im} \omega t(y, x)) \tag{23}
\end{equation*}
$$

This inequality agrees with (7) and (8). For a singular function of type (14) defined in an open set $Y \subset X$ the estimate

$$
\int_{X} h \exp (2 \operatorname{Im} \omega t(y, x))\left(\left|\hat{A}^{(1)}(x ; \omega)\right|^{2}+\left|\dot{\omega} \square_{\omega} A^{(1)}(x ; \omega)\right|^{2}\right) d x \leq \frac{C}{|\omega|}
$$

holds for an arbitrary continuous function $h$ compactly supported in $Y$. It can be checked by the method of Lemma 11.1 below. In the same way we can estimate the integrals of $\left|\nabla B_{j k}\right|^{2}$ in (21), since each kernel $\nabla B_{j k}$ is of type (14) or of type (11). Summing these inequalities over $\alpha$ and taking into account (21), we get

$$
\begin{equation*}
\int_{X} h \exp (2 \operatorname{Im} \omega t(y, x))\left(\left|\widehat{P}_{D}(y, x ; \omega)\right|^{2}+\left|\omega \widehat{Q}_{D}(y \cdot x ; \omega)\right|^{2}\right) d x \leq \frac{C}{|\omega|} \tag{24}
\end{equation*}
$$

for any continuous $h$ compactly supported in $D_{0}$.

## 8. Phase functions for constant velocity

Step 6. We extend the above construction for the domain $D \times X$. First, we choose a special generating function for $\Lambda_{\infty} \doteq \Lambda \cap \pi^{-1}(D \times X \backslash D)$. Denote by $\Gamma_{s}$ the solution of (6) with the initial data $x_{s}=y$ and $\xi_{s}=(\mathbf{n}(y) \cos s, \mathbf{n}(y) \sin s), 0 \leq s<2 \pi$. By (*) $\mathbf{n}=\mathbf{n}_{0}$ is constant in $X \backslash D$; the Hamiltonian system has the form

$$
\frac{d x}{d r}=-\xi, \quad \frac{d t}{d r}=\mathbf{n}^{2} . \quad \frac{d \xi}{d r}=0, \quad \tau= \pm 1
$$

and the ray $\gamma(s) \doteq \pi_{y}\left(\Gamma_{s}\right)$ is a straight line in $X \backslash D$. We assume that $D_{0}$ is a convex neighbourhood of $D$ with smooth boundary $\partial D_{0}$. Denote by $\left(x_{0}(s), \xi_{0}(s)\right)$ the point
in $\Gamma_{s}$, where the ray $\gamma(s)$ reaches the boundary. The vector $\xi(s)=-d x / d r$ cannot be tangent to the convex curve $\partial D_{0}$, because otherwise the ray $\gamma(s)$ could not enter $D$. The conditions (*) and (**) imply that $x_{0}$ and $\xi_{0}$ are $C^{2}$-functions in $D \times S^{1}$. The ray $\gamma(s) \backslash D_{0}$ is parameterized by

$$
\begin{equation*}
x(s, r)=-\xi_{0}(s) r+x_{0}(s) \text { and } t(s, r)=\mathbf{n}^{2} r+t_{0}(s) \quad r \geq 0 \tag{25}
\end{equation*}
$$

hence the parameter $r$ vanishes just at $\partial D_{0}$. The ray $\gamma_{;}(s)$ and its direction $\xi_{0}(s)$ will be called critical if $\xi_{0}^{\prime}(s)=0$. The union $\Xi$ of all critical directions has zero angular measure. The caustic set, i.e. the image of the critical set of $\pi_{y}$. is contained in $V \cup G$, where $V$ is an arbitrary open conic neighbourhood of $\Xi$ and $G$ is a compact set.

Consider the phase function $\phi(x, t: s)=t-\varphi(x: s)$. where

$$
\begin{equation*}
\varphi(x ; s) \doteq t_{0}(s)+\left\langle\xi_{0}(s) \cdot x_{0}(s)-x\right\rangle=t(s . r)+\left\langle\xi_{0}(s) \cdot x(s, r)-x\right\rangle \tag{26}
\end{equation*}
$$

and $\langle\cdot, \cdot\rangle$ stands for the inner product in $X$. This function satisfies the eikonal equation $\left|\nabla_{x} \varphi\right|^{2}=\left|\xi_{0}\right|^{2}=\mathbf{n}^{2}$. The function $o$ vanishes in $\Lambda_{\times}$for $t=\tau(s, r)$ and

$$
d \varphi(x ; s)=\left[d t_{0}(s)+\left\langle\xi_{0}(s) \cdot d x_{0}(s)\right\rangle\right]-\left\langle\xi_{0}(s) \cdot d x\right\rangle+\left\langle d \xi_{0}(s) \cdot x_{0}(s)-x\right\rangle .
$$

The first bracket vanishes in the Lagrange manifold $\Lambda_{x}$. We have $\varphi_{s}^{\prime}(x ; s)=$ $\left\langle\xi_{0}^{\prime}(s), x_{0}(s)-x\right\rangle$ and by the eikonal equation. $\left\langle\xi_{0}^{\prime} \cdot \xi_{0}\right\rangle=0$. Consequently the equation

$$
-\phi_{s}^{\prime}=\varphi_{s}^{\prime}=\left\langle\xi_{0}^{\prime}(s), x_{0}(s)-x(s, r)\right\rangle=r\left\langle\xi_{0}^{\prime}(s), \xi_{0}(s)\right\rangle=0
$$

holds in $\Lambda_{\infty}$. Vice versa, for any noncritical direction. the equation $\psi_{s}^{\prime}=0 \mathrm{im}$ plies that $x_{0}(s)-x=r \xi_{0}(s)$ for some $r \geq 0$ and $o(x . t: s)=t-\tau(s . r)$. Moreover, $d_{x} o=$ $-d_{x} \varphi=\left\langle\xi_{0}, d x\right\rangle$; hence $\phi$ generates the noncritical part of $\Lambda_{x}$.

Set $\sigma \doteq \xi_{0} \times \xi_{0}^{\prime}$; then $|\sigma| \doteq \mathbf{n}\left|\xi_{0}^{\prime}\right|$. For a vector $\tau=\left(v_{1}, v_{2}\right)$ we let $v^{*}=\left(-v_{2}, v_{1}\right)$ : then $\left\langle v^{*}, u\right\rangle=v \times u$, where $u$ is an arbitrary vector $u$ and $\times$ denotes the cross product.

Lemma 8.1. The quantity $j(s, r)=r \sigma+x_{0}^{\prime} \times \xi_{0}$ is equal to the divergence of the family of rays and

$$
\begin{equation*}
s_{x}^{\prime}=\frac{\xi_{0}^{*}}{j}, \quad r_{x}^{\prime}=\frac{1}{j}\left(-r \xi_{0}^{\prime}+x_{0}^{\prime}\right)^{*} . \tag{27}
\end{equation*}
$$

Proof. The divergence equals the Jacobian

$$
\begin{equation*}
j(s, r)=\operatorname{det} \frac{\partial x}{\partial(r . s)}=x_{r}^{\prime} \times x_{s}^{\prime}=-\xi_{0} \times\left(-r \xi_{0}^{\prime}+x_{0}^{\prime}\right)=r \sigma+x_{0}^{\prime} \times \xi_{0} \tag{28}
\end{equation*}
$$

It follows that there is focusing near any critical ray. Calculate the transport system in ( $s, r$ )-coordinates for the phase function (26). We have $\nabla_{x \varphi}=-\xi_{0}(s)$ and

$$
\Delta \varphi=\left\langle\nabla, \nabla_{x} \psi\right\rangle=-\left\langle\xi_{0}^{\prime} \cdot s_{x}^{\prime}\right\rangle=-\frac{1}{j}\left\langle\xi_{0}^{\prime} \cdot \xi_{0}^{*}\right\rangle=\frac{\sigma}{j} .
$$

Hence (12) appears as

$$
\begin{equation*}
\partial_{r}\left(a|j|^{1 / 2}\right)=0 \tag{29}
\end{equation*}
$$

It follows that $L_{r}\left(|a|^{2} d x\right)=0$, where $L_{r}$ is the Lie derivative with respect to the field $\partial_{r}=\partial / \partial r$.

Lemma 8.2. The equations $x_{0}^{\prime} \times \xi_{0}=0$ and $\xi_{0}^{\prime}=0$ have no common solution.
Proof. The vector $x_{0}^{\prime}$ is tangent to $\partial D_{0}$ and the vector $\xi_{0}$ is not. Therefore the first equation is equivalent to $x_{0}^{\prime}=0$. This together with the second equation implies that $s$ is no more a local coordinate in $\Lambda$ for a point $\left(s_{0} \cdot r_{0}\right)$. which contradicts Liouville's theorem.

## 9. The parametrix in the nonperturbed domain

Step 7. Choose small numbers $\varepsilon>0$ and $\delta>0$ such that

$$
\begin{equation*}
|\sigma| \leq 3 \varepsilon \quad \text { implies } \quad\left|\xi_{0} \times x_{0}^{\prime}\right| \geq \delta \tag{30}
\end{equation*}
$$

and set $\mu \doteq \max _{0 \leq s \leq 2 \pi}\left|\xi_{0}(s) \times x_{0}^{\prime}(s)\right|>0$. We can assume that the set $D_{0}$ chosen in Step 4 contains a $2 \mu \mathbf{n} \varepsilon^{-1}$-neighbourhood of $D$ where $\mathbf{n}=\mathbf{n}_{0}$. Take the covering of $\Lambda_{\infty}$ by the sets

$$
\begin{aligned}
& \Lambda_{>} \doteq\{(s, r):|\sigma(s)| \geq \varepsilon\} \\
& \Lambda_{<} \doteq\left\{(s, r):|\sigma(s)| \leq 3 \varepsilon \text { and either } r|\sigma(s)| \leq \frac{1}{2} \delta \text { or } r|\sigma(s)| \geq 2 \mu\right\} \\
& \Lambda_{b} \doteq\left\{(s, r):|\sigma(s)| \leq 3 \varepsilon \text { and } \frac{1}{3} \delta \leq r|\sigma(s)| \leq 3 \mu\right\} .
\end{aligned}
$$

By (25), $\mathbf{n} r=\left|x(s, r)-x_{0}(s)\right| \geq 2 \mu \mathbf{n} \varepsilon^{-1}$ in $X \backslash D_{0}$. which yields $r \geq 2 \mu \varepsilon^{-1}$. Therefore $r|\sigma| \geq 2 \mu \geq\left|\xi_{0} \times x_{0}^{\prime}\right|+\mu$ and $|j| \geq r|\sigma|-\left|x_{0}^{\prime} \times \xi_{0}\right| \geq \mu$ in $\Lambda_{>}$. Hence there is no focusing in $\Lambda_{>}$and we can use $x$ as a local coordinate system. We have by (25),

$$
\begin{equation*}
\mathbf{n} r-d_{0} \leq|x-y| \leq \mathbf{n}_{0} r+d_{0}, \quad d_{0} \doteq \max _{x \in \partial D_{0}}|x-y| \tag{31}
\end{equation*}
$$

since $\left|\xi_{0}\right|=\mathbf{n}$. Take a smooth function $h_{>}$in $\mathbf{R}$ such that $h_{>}(\tau)=1$, if $|\tau| \geq 3 \varepsilon$, and $h_{>}(\tau)=0$, if $|\tau| \leq 2 \varepsilon$. Set

$$
\begin{equation*}
A_{>}(x, t) \doteq \sum_{x(s, r)=x}(t-\tau(s, r))_{ \pm}^{-1 / 2} e(t-\tau(s, r)) a(s, r) h_{>}(\sigma(s)), \quad \pm 1=\operatorname{sgn} j(s, r) \tag{32}
\end{equation*}
$$

where the sum is taken over all points $(s, r) \in \Lambda_{>}$such that $x(s, r) \doteq-\xi_{0}(s) r+x_{0}(s)=$ $x$, and the amplitude function $a_{0}$ is smooth and satisfies (29). We take the initial data for $a$ from the consistency condition with the kernel $P_{D}$. This means that the difference $A_{>}-P_{D}$ equals the sum of terms of type (17) in a neighbourhood of $\partial D_{0}$.

Lemma 9.1. The number $N(x)$ of terms in (32) that are nonvanishing for $x \in X \backslash D_{0}$ is uniformly bounded.

Proof. We call an interval $I \subset S^{1}$ univalent if $|\sigma(s)| \geq \varepsilon$ for $s \in I$. The inequality $|j(s, r)|>0$ holds for $r \geq \mu \varepsilon^{-1}$ in an arbitrary univalent interval. Consider the mapping

$$
\begin{equation*}
\pi_{y}: I \times\left[\mu_{0}, \infty\right) \longrightarrow X . \quad(s, r) \longmapsto x(s, r) \tag{33}
\end{equation*}
$$

which is affine with respect to $r$. It is an imbedding, since the Jacobian does not vanish and $I$ is connected. Hence for any $x \in X \backslash D_{0}$ there is no more than one solution to the equation $x=x(s, r)$. If the whole circle $S^{1}$ is univalent, the number of terms is equal to 1 . Suppose the opposite, i.e. there is a point $s^{\prime}$ such that $\left|\sigma\left(s^{\prime}\right)\right|<\varepsilon$. We call an interval $I=\left(s_{0}, s_{1}\right)$ good, if it is univalent and $\left|\sigma\left(s_{0}\right)\right|=$ $\left|\sigma\left(s_{1}\right)\right|=2 \varepsilon$. We have $|\sigma(s)|<2 \varepsilon$, if a point $s$ does not belong to a good interval, and the corresponding term in (32) vanishes since $h_{>}(\sigma)=0$. If two good intervals have a common point, then their union is again a good interval and the number of nonvanishing terms in (32) is bounded by the number $N$ of maximal good intervals. We show that $N<\infty$. Indeed, for arbitrary consecutive maximal good intervals $\left(s_{0}, s_{1}\right),\left(s_{2}, s_{3}\right) \subset S^{1}$ such that $s_{1}<s_{2}$, there is a point $\tau \in\left(s_{1}, s_{2}\right)$ such that $|\sigma(\tau)|<\varepsilon$, whereas $\left|\sigma\left(s_{1}\right)\right|=\left|\sigma\left(s_{2}\right)\right|=2 \varepsilon$. These relations imply that $s_{2}-s_{1}>2 \varepsilon / \max \left|\sigma^{\prime}(s)\right|$. The same inequality holds for $s_{0}+2 \pi-s_{3}$. Therefore $N<\pi \varepsilon^{-1} \max \left|\sigma^{\prime}\right|$.

We have the inequality

$$
\int_{\Gamma(\tau)} \sum_{x(s, r)=x}\left|a(s, r) h_{>}(\sigma)\right|^{2} d l \leq C, \quad \Gamma(\tau) \doteq\{x: t(y, x)=\tau\}
$$

where $d l$ denotes the arc element of $\Gamma(\tau)$ and $C$ does not depend on $\tau \geq \tau_{0}$. This follows from the equation $d l=\left(1+O\left(r^{-1}\right)\right) d x / d r$, where the form $d r$ and the density
$|a|^{2} d x$ are constant along any ray by virtue of (29). The number of terms in the sum is uniformly bounded by Lemma 9.1. By (23) and (32) we find

$$
\begin{equation*}
\int_{\Gamma(\tau)}\left|\hat{A}_{>}(x ; \omega)\right|^{2} d l \leq \frac{C}{|\omega|} \int_{\Gamma(\tau)}\left|a(s, r) h_{>}(\sigma)\right|^{2} d l \leq \frac{C}{|\omega|} \tag{34}
\end{equation*}
$$

and

$$
\begin{equation*}
\int_{\Gamma(\tau)}\left|\hat{A}_{>}(x ; \omega)\right|^{2} d l \leq \frac{C}{|\dot{\omega}|} \tag{35}
\end{equation*}
$$

where the constant $C$ does not depend on $\tau$ and $\omega$.
Step 8. Choose a smooth even function $h_{<}$in $\mathbf{R}$ such that $h_{<}=1$ in $\left[0, \frac{1}{3} \delta\right] \cup$ $[\mu+2, \infty)$ and $h_{<}=0$ in $\left[\frac{1}{2} \delta, \mu+1\right]$. Define the singular function which is similar to (32):

$$
\begin{equation*}
A_{<}(x, t)=\sum_{x(s, r)=x}(t-t(s, r))_{ \pm}^{-1 / 2} e(t-\tau(s . r)) a(s, r)\left(1-h_{>}(\sigma(s))\right) h_{<}(r \sigma(s)) \tag{36}
\end{equation*}
$$

where $\pm 1=\operatorname{sgn} j(s, r)$ and the amplitude $a$ is defined as in Step 7. Note that $j \neq 0$ in $\operatorname{supp}\left(1-h_{>}\right) h_{<}$; hence the equation $x(s . r)=x$ has a locally smooth solution $s=s(x)$, $r=r(x)$. The function $A_{<}$is continued to $D \times\left(X \backslash D_{0} \times \mathbf{R}\right)$ by setting $A_{<}=0$ for $x \in X \backslash\left(D_{0} \cup \pi_{y}\left(\Lambda_{<}\right)\right)$. It is smooth in $D \times X$.

Lemma 9.2. The number of nonzero terms in (36) is uniformly bounded for $x \in X \backslash D_{0}$.

Proof. We argue similarly to the proof of Lemma 9.1. We call an interval $J \subset S^{1}$ critical if $|\sigma| \leq 3 \varepsilon$ in $J$. Let $J$ be a critical interval. If $s \in J$, the inequality $j(s, r) \geq \frac{1}{2} \delta$ holds for $r \leq \delta / 2|\sigma|$, because of (30), whereas $j(s, r) \geq 1$ for $r \geq 2 \mu|\sigma|^{-1}$, $\sigma \neq 0$. Therefore the mapping $J \times\left[0, \delta(2|\sigma|)^{-1}\right] \cup\left[2 \mu|\sigma|^{-1} . \infty\right) \rightarrow X$ like (33) is an imbedding and for any $x \in X$ there is no more than two solutions to $x(s, r)=x$, if $s \in J,(s, r) \in \Lambda_{<}$. If $|\sigma| \geq 2 \varepsilon$ everywhere in $S^{1}$, the sum (36) vanishes since $h_{>}(\sigma)=$ 1. Suppose the opposite. We call a critical interval $J=\left(s_{0}, s_{1}\right)$ good, if $\left|\sigma\left(s_{0}\right)\right|=$ $\left|\sigma\left(s_{1}\right)\right|=2 \varepsilon$. The number of nonzero terms is bounded by the number $N$ of maximal good intervals. We show that this number is finite. Two different maximal good intervals are disjoint. If a point $s$ does not belong to a good interval, we have $|\sigma(s)|>3 \varepsilon$ and the corresponding term in (36) vanishes since $h_{>}(\sigma)=1$. For arbitrary maximal good intervals $\left(s_{0}, s_{1}\right)$ and $\left(s_{2}, s_{3}\right)$, there is a point $\tau \in\left(s_{1}, s_{2}\right)$ such that $|\sigma(\tau)|>3 \varepsilon$, whereas $\left|\sigma\left(x_{1}\right)\right|=\left|\sigma\left(x_{2}\right)\right|=2 \varepsilon$. This yields $s_{2}-s_{1} \geq 2 \varepsilon\left(\max \left|\sigma^{\prime}(s)\right|\right)^{-1}$ and $N \leq \pi \varepsilon^{-1} \max \left|\sigma^{\prime}\right|$.

## 10. The parametrix near caustics

Step 9. The integral

$$
\begin{equation*}
A_{b}(x, t) \doteq \int_{\mathbf{R}} e(\phi(x, t ; s)) \operatorname{Re}\left(p_{0}\left(\phi_{,}(x, t ; s)\right) a_{\supset}(x, s)\right) h_{\supset}(s, r) d s \tag{37}
\end{equation*}
$$

represents the parametrix in the focusing area $\Lambda_{3}$. Here $e$ is a function as in Lemma 5.2,

$$
\begin{aligned}
\phi_{b}(x, t ; s) & =t-\psi_{j}(x, s) . \\
\psi_{y}(x, s) & =t_{0}(s)+\left\langle\xi_{0}(s) \cdot x_{0}(s)-x\right\rangle \\
h_{\partial}(s, r) & \doteq\left(1-h_{>}(\sigma(s))\right)\left(1-h_{<}(\sigma(s) r)\right) . \\
r & \doteq \frac{1}{\mathbf{n}^{2}}\left\langle\xi_{0}(s), x_{0}(s)-x\right\rangle
\end{aligned}
$$

and the amplitude $a_{0}$ is defined below. We have $\left(\mathbf{n}_{0}\right)_{r}^{\prime}=0$ and $\varphi_{s}^{\prime}=0$ for $x=x(s, r)$. Therefore, since $\Delta \varphi_{b}=0$, (15) looks as follows:

$$
\begin{equation*}
2 \partial_{r} a_{\nu}=2\left\langle\xi_{0} . \nabla a_{\jmath}\right\rangle=\left(\hat{\gamma}_{j}\right)_{s}^{\prime} b . \tag{38}
\end{equation*}
$$

Introduce the variable $u=\mathbf{n}^{-2} \xi_{0} \times\left(x-x_{0}\right)$; we have $\left(\hat{\gamma^{\prime}}\right)_{s}^{\prime}=\left\langle\xi_{0}^{\prime} \cdot x-x_{0}\right\rangle=\sigma u$. The function $a_{b}$ should satisfy (38) in the surface $u=0$. We extend this function by setting $a_{b}(x, s, r) \doteq a_{b}(x(s, r), s) g(u), x(s . r)=x_{0}(s)-r \xi_{0}(s)$. where $g \in \mathcal{D}((-1.1))$ and $g=1$ in a neighbourhood of the origin. It is easy to see that (38) now is satisfied with $b=0$ for any $u$.

Next we find the consistency condition for amplitudes in the set $\pi\left(\Lambda_{-}\right)$, where $\Lambda_{\sharp} \doteq\left(\Lambda_{>} \cup \Lambda_{<}\right) \cap \Lambda_{b}$. First we change the integration variable $s$ in (37) to $u=u(x, s)$. From $u_{s}^{\prime}=\mathbf{n}^{-2}\left(\xi_{0}^{\prime} \times\left(x-x_{0}\right)-\xi_{0} \times x_{0}^{\prime}\right)$, by substituting $x=x_{0}-r \xi_{0}+u \xi_{0}^{*}$, we obtain $u_{s}^{\prime}=\mathbf{n}^{-2}\left(r \sigma-\xi_{0} \times x_{0}^{\prime}\right)=\mathbf{n}^{-2} j(s, r)$. We have $\frac{1}{2} \delta \leq r|\sigma| \leq \frac{2}{3} \delta$ or $1 \leq r|\sigma| \leq 2 \mu+2$ in $\Lambda_{r}$. Therefore the derivative $v_{s}^{\prime}$ is bounded from above and below, the function $s=s(x, v)$ is smooth, and we can write

$$
A_{b}(x, t) \doteq \int_{\mathbf{R}} e\left(\phi_{\partial}(x, t ; s)\right) \operatorname{Re}\left(p_{0}(o,(x, t: s)) a_{\gamma}(x . s)\right) \frac{h_{,}(s, r)}{|j(s . r)|} \mathbf{n}^{2} d u
$$

Next we calculate this integral as in Step 4. We have $\partial_{1: \hat{\gamma}}=\partial_{s} \hat{\varphi}_{y}\left(u_{s}^{\prime}\right)^{-1}=\mathbf{n}^{2} \sigma u j^{-1}$; hence $\varphi_{b}(x, s)=\varphi(x)+\chi(s, r) u^{2}$. $\chi \doteq \mathbf{n}^{2} \sigma(2 j)^{-1}$. where $\varphi(x)=\tau(s, r)$ is the phase function of $A_{>}$. The function $|\chi|$ is smooth and bounded from below since $|j| \geq \frac{1}{2} \delta$ and $|\sigma| \geq \delta / 2 r$. As $\phi_{\nu}(x, t ; s) \equiv t-\varphi_{\nu}(x ; s)$, we can integrate by method (19), getting

$$
\begin{align*}
A_{b}(x, t) & =\int_{\mathbf{R}} e\left(\phi+\chi v^{2}\right) \operatorname{Re}\left(p_{0}\left(\phi+\chi v^{2}\right) a_{,}(x, s)\right) \frac{h_{,}(x)}{|j(s . r)|} \mathbf{n}^{2} d u \\
& =\operatorname{Re}\left(\left(\phi_{+}^{-1 / 2}+\varepsilon \imath \phi_{-}^{-1 / 2}\right) b(x, 0) h_{,}(x)\right)-\int_{\mathbf{R}} \operatorname{Re} p_{1}\left(\phi_{,}\right) c_{v}^{\prime}(x, v) h_{,}(x) d u \tag{39}
\end{align*}
$$

where $\phi=t-\varphi(x), \varepsilon=-\operatorname{sgn} \chi$ and

$$
b(x, v) \doteq 2^{1 / 2} \pi \mathbf{n} a_{b}(x, s)|j \sigma|^{-1 / 2}=b(x, 0)+2 v c(x, v)
$$

Now we can write the consistency condition

$$
\begin{equation*}
a(s, r)=b(x, 0)=2^{1 / 2} \pi \mathbf{n} a,(x . s(x .0))|j \sigma|^{-1 / 2} \tag{40}
\end{equation*}
$$

where $x=x(s, r), u=0$ and $a$ is the amplitude of $A_{<}$. We see that the product $a(s, r)|j|^{1 / 2}=2^{1 / 2} \pi \mathbf{n} a_{b}(x, s(x, 0))|\sigma|^{-1 / 2}$ does not depend on $r$, since $a$, satisfies (38).

## 11. The energy in the caustic area

Now we estimate the integral of $\hat{A}_{y}$. First calculate the Fourier transform of the functions $p_{k}(\phi)$, where $\phi=t-\varphi$. We have for $\omega>0$.

$$
\begin{aligned}
& F\left(p_{0}(\phi)\right)=-2 \pi \imath \exp (-\imath \iota \imath) . \\
& F\left(p_{1}(\phi)\right)=-\frac{2 \pi \imath}{\omega} \exp (-\imath \omega \vartheta) . \\
& F\left(\bar{p}_{0}(\phi)\right)=F\left(\bar{p}_{1}(o)\right)=0 ;
\end{aligned}
$$

for $\omega<0$ we need only to replace $p_{k}$ by $\bar{p}_{k}$. Therefore

$$
\hat{A}_{b}(x ; \omega)=\int_{\mathbf{R}} \exp \left(-\omega^{\prime} \varphi_{\rangle}(x, s)\right) \tilde{a}(x, s) d s+R_{\nu}(x ; \omega)
$$

where $\tilde{a}(x, s)=2 \pi i a_{b}(x, s) h_{\partial}(s, r)$ and $\left|R_{j}(x ; \omega)\right| \leq C \mid \omega^{-1} \exp (-\operatorname{Im} \omega t(y, x))$. We get

$$
\begin{aligned}
\left|\hat{A}_{b}(x ; \omega)\right|^{2} \leq & \int_{\mathbf{R}} \int_{\mathbf{R}} \exp \left(-\imath \operatorname{Re} \omega \Phi\left(x ; s, s^{\prime}\right)-\operatorname{Im} \omega \Psi\left(x ; s, s^{\prime}\right)\right) a\left(x ; s, s^{\prime}\right) d s d s^{\prime} \\
& +\left|R_{b}(x, \omega)\right|^{2}
\end{aligned}
$$

where

$$
\begin{aligned}
& \Phi\left(x ; s, s^{\prime}\right) \doteq \varphi_{b}(x, s)-\varphi_{b}\left(x, s^{\prime}\right)=t_{0}(s)-t_{0}\left(s^{\prime}\right)+\xi_{0}(s)\left(x_{0}(s)-x\right)-\xi_{0}\left(s^{\prime}\right)\left(x_{0}\left(s^{\prime}\right)-x\right) \\
& \Psi\left(x ; s, s^{\prime}\right) \doteq \varphi_{b}(x, s)+\varphi_{b}\left(x, s^{\prime}\right) \\
& a\left(x ; s, s^{\prime}\right) \doteq \tilde{a}(x, s) \overline{\bar{a}}\left(x, s^{\prime}\right)
\end{aligned}
$$

Integrating over $\Gamma(\tau)$ for $\tau \geq \tau_{0}$, yields

$$
\int_{\Gamma(\tau)}\left|\hat{A}_{b}(x ; \omega)\right|^{2} d l \leq I(\tau \cdot \omega)+R(\tau \cdot \dot{\omega})
$$

where

$$
\begin{aligned}
I(\tau, \omega) & \doteq \int_{\Gamma(\tau)} \int_{\mathbf{R}} \int_{\mathbf{R}} \exp \left(-\imath \operatorname{Re} \omega \Phi\left(x ; s, s^{\prime}\right)-\operatorname{Im} \omega \Psi\left(x: s . s^{\prime}\right)\right) a\left(x: s . s^{\prime}\right) d s d s^{\prime} d l \\
R(\tau, \omega) & \doteq \int_{\Gamma(\tau)}\left|R_{\zeta}(x, \omega)\right|^{2} d l
\end{aligned}
$$

## Lemma 11.1. The following inequalities hold

$$
\begin{equation*}
|I(\tau, \omega)| \leq \frac{C}{|\omega|} \tau \exp (-\operatorname{Im} \omega \tau) . \quad R(\tau, \dot{\prime}) \leq \frac{C}{|\dot{\omega}|^{2}} \tau \exp (-\operatorname{Im} \omega \tau), \quad \tau \geq 0 \tag{41}
\end{equation*}
$$

Proof. We will apply the stationary phase evaluation. The function $\varphi_{j}$ belongs to $C^{2}$ since $\xi_{0}, x_{0}$ are of this class (see Step 9). Therefore the phase function $\Phi$ is contained in $C^{2}$ too. Find the critical set of $\left\{\Phi_{l}=\Phi_{s}=\Phi_{s^{\prime}}=0\right\}$ in the manifold $\Gamma(\tau) \times S^{1} \times S^{1}$. The $l$-derivative can be written as $\partial_{l} \doteq|x-y|^{-1}(x-y) \times \partial_{x}$. Therefore the first equation is equivalent to $\left(\xi_{0}(s)-\xi_{0}\left(s^{\prime}\right)\right) \times(x-y)=0$ and implies $\xi_{0}(s)=$ $\xi_{0}\left(s^{\prime}\right)$. The second and third equations give $\left\langle\xi_{0}^{\prime}(s) \cdot x-x_{0}(s)\right\rangle=\left\langle\xi_{0}^{\prime}(s) \cdot x-x_{0}\left(s^{\prime}\right)\right\rangle=0$ which imply $x=x(s, r)$ for some $r \geq 0$ and $x_{0}(s)=x_{0}\left(s^{\prime}\right)$. It follows that $\Gamma(s)=\Gamma\left(s^{\prime}\right)$, hence $s^{\prime}=s$. Calculating the second derivatives in $s$ and $l$ on the diagonal $\Delta(\tau)$ in $\Gamma(\tau) \times S^{1} \times S^{1}$ defined by $s=s^{\prime}$, we get

$$
\begin{align*}
\Phi_{l s}^{\prime \prime} & =\xi_{0}^{\prime}(s) \times|x-y|^{-1}(x-y)=\left(-\frac{1}{\mathbf{n}}+O\left(\frac{1}{\tau}\right)\right) \xi_{0}^{\prime}(s) \times \xi_{0}(s)=\left(\frac{1}{\mathbf{n}}+O\left(\frac{1}{\tau}\right)\right) \sigma,  \tag{42}\\
\Phi_{s s}^{\prime \prime} & =r\left\langle\xi_{0}^{\prime \prime}, \xi_{0}\right\rangle+\left\langle\xi_{0}^{\prime}, x_{0}^{\prime}\right\rangle=-r\left|\xi_{0}^{\prime}\right|^{2}+\left\langle\xi_{0}^{\prime} \cdot x_{0}^{\prime}\right\rangle=\frac{-r \sigma^{2}+\sigma \xi_{0} \times x_{0}^{\prime}}{\mathbf{n}^{2}}=\frac{\sigma j}{\mathbf{n}^{2}} \\
\Phi_{l l}^{\prime \prime} & =0 \\
\operatorname{det} \nabla_{l s}^{2} \Phi & =\left(-1+O\left(\frac{1}{\tau}\right)\right) \frac{\sigma^{2}}{\mathbf{n}^{2}}
\end{align*}
$$

since $r^{-1} x \rightarrow \xi_{0}$, as $\tau \rightarrow \infty$, and $r / \tau=1+O\left(\tau^{-1}\right)$. We have $\left|\Phi_{s s}^{\prime \prime}\right| \leq 4 \mu|\sigma| \mathbf{n}^{-2}$, since $|j| \leq 4 \mu$ in $\Lambda_{b}$; hence all second derivatives in (42) are equal to $O(|\sigma|)$. We apply Proposition 17.1 to the function $\Phi$ for fixed $s$ and $x \in \Gamma(\tau)$. There is only one critical point $x_{*}=x_{0}(s)-r \xi_{0}(s)$. We have $\left\|b^{-1}\right\|=O\left(|\sigma|^{-1}\right)$ for $b=\frac{1}{2} d^{2} \Phi(x ; s, s)$. From (42) we see that $\left|\nabla_{l s}^{3} \Phi\right|=O(1)$. Therefore the set $\left\{(s . x):\left|s-s^{\prime}\right|+\left|x-x_{*}\right| \leq c\left|\sigma\left(s^{\prime}\right)\right|\right\}$ is contained in $U_{0}$. Then we apply Proposition 17.2. All derivatives of first and second order of $a$ are bounded in $\Lambda_{\nu}$. We have $\hat{y}^{2}(x, s)>t(y . x)$, since the caustic is always behind the first front of a wave. This yields $\Psi\left(x: s . s^{\prime}\right)>2 t(y . x)$. Therefore by Proposition 17.3 we obtain $|I(\tau, \omega)| \leq C|\omega \sigma|^{-1} \exp (-\operatorname{Im} \omega \tau)$. Taking into account that $\frac{1}{3} \delta \leq r|\sigma|$, we estimate $|\sigma|^{-1} \leq C r \leq C \tau$. This proves the estimate (41) for $I(\tau, \omega)$. The estimate for the remainder $R$ is more simple.

## 12. The parametrix in the homogeneous domain

Step 10. We set $P_{\infty}=A_{>}+A_{<}+A_{,}$, where we choose the amplitude functions in (32), (36) and (37) according to the rule (40). It follows from (34), (35) and (41)
that

$$
\begin{equation*}
\int_{\Gamma(\tau)}\left|\widehat{P}_{\infty}(x ; \omega)\right|^{2} d l \leq \frac{C}{|\omega|} \tau \exp (-\operatorname{Im} \omega \tau) \tag{43}
\end{equation*}
$$

We estimate $Q_{\infty} \doteq \square P_{\infty}$ in $D \times\left(X \backslash D_{0} \times \mathbf{R}\right)$. We have $Q_{\infty}=\square\left(A_{>}+A_{<}\right)+\square A_{b}$ and (44)

$$
\begin{aligned}
\square\left(A_{>}(x, t)+A_{<}(x, t)\right) \doteq & \sum_{x(s, r)=x} e(\phi) \phi_{ \pm}^{1 / 2} \Delta a\left(1-h_{\gamma}\right) \\
& -2 \sum_{x(s, r)=x}\left\langle\nabla\left(e(\phi) \phi_{ \pm}^{1 / 2} a\right), \nabla h_{\gamma}\right\rangle-\sum_{x(s . r)=x} e(\phi) \phi_{ \pm}^{1 / 2} a \Delta h_{\supset}
\end{aligned}
$$

where $a$ is defined in (32) and

$$
\begin{align*}
\square A_{b}(x, t)= & \int_{\mathbf{R}} e(\phi) \operatorname{Re}\left(p_{1}(\phi)(\Delta+\mathbf{m}) a_{\jmath}\right) h_{\nu} d s  \tag{45}\\
& +2 \int_{\mathbf{R}}\left(\nabla\left(\operatorname{Re} p_{1}(\phi) a_{\jmath}\right) \cdot \nabla h_{\supset}\right\rangle d s+\int_{\mathbf{R}} \operatorname{Re} p_{1}(\phi) a_{\supset} \Delta h_{,} d s
\end{align*}
$$

where $a_{b}$ is defined in (38). We want to estimate the Fourier transform of the six terms in (44) and (45). Write the sum as $Q_{\infty}=Q^{0}+Q^{1}+Q^{2}$, where

$$
\begin{aligned}
& Q^{0} \doteq \sum_{x(s, r)=x} e(\phi) \phi_{ \pm}^{1 / 2}\left(1-h_{b}\right)(\Delta+\mathbf{m}) a+\int_{\mathbf{R}} e(\phi) \operatorname{Re}\left(p_{1}(\phi)(\Delta+\mathbf{m}) a_{\rho}\right) h_{,} d s, \\
& Q^{1} \doteq-2 \sum_{x(s, r)=x}\left\langle\nabla\left(e(\phi) \phi_{ \pm}^{1 / 2} a\right), \nabla h_{\jmath}\right\rangle+2 \int_{\mathbf{R}}\left\langle\nabla e(o) \operatorname{Re}\left(p_{1}(\rho) a_{\jmath}\right) . \nabla h_{\jmath}\right\rangle d s, \\
& Q^{2} \doteq-\sum_{x(s, r)=x} e(\phi) \phi_{ \pm}^{1 / 2} a \Delta h_{\mathrm{\jmath}}+\int_{\mathbf{R}} \operatorname{Re} p_{1}(\phi) a, \Delta h_{,} d s .
\end{aligned}
$$

First, check the inequality

$$
\begin{equation*}
\int_{\Gamma(\tau)}\left|\widehat{Q}^{0}(x ; \omega)\right|^{2} d l \leq \frac{C}{|\omega|^{3}} \tau \exp (-\operatorname{Im} \omega \tau) \tag{46}
\end{equation*}
$$

For the first term in $Q^{0}$ this follows from (23). We estimate the second term similarly to (41), taking into account that derivatives of the amplitude $a$, are bounded.

The terms $Q^{1}$ and $Q^{2}$ contain derivatives of the cutting functions $h_{>}$and $h_{<}$. The derivatives $\nabla h_{>}(\sigma)$ and $\Delta h_{>}(\sigma)$ are bounded by (27) and (31) since
the dominator $j$ has a positive lower bound. By similar arguments we show that $\nabla h_{<}(\sigma r)=O(|x|)$ and $\Delta h_{<}(\sigma r)=O\left(|x|^{2}\right)$. as $|x| \rightarrow x$ : hence also

$$
\begin{equation*}
h_{\mathrm{b}}=O(1), \quad \nabla h,=O(|x|) \quad \text { and } \quad \Delta h_{\imath}=O\left(|x|^{2}\right) \tag{47}
\end{equation*}
$$

Therefore

$$
\begin{equation*}
\int_{\Gamma(\tau)}\left|\widehat{Q}^{2}(x ; \omega)\right|^{2} d l \leq \frac{C}{|\dot{\omega}|^{3}} \exp (-\tau \operatorname{Im} \omega)\left(\tau^{5}+\tau\right) \tag{48}
\end{equation*}
$$

The function $Q^{1}$ is supported in $\pi\left(\Lambda_{2}\right)$, since $\nabla h_{3}=0$ in the complement. We write this function in a different form where no kernel $\nabla O_{ \pm}^{-1 / 2}=c \nabla(\phi) \phi_{ \pm}^{-3 / 2}$ appears. For this we apply (39):

$$
\int_{\mathbf{R}} e(\phi) \operatorname{Re}\left(p_{0}\left(\phi_{j}\right) a_{j}\right) d s=\operatorname{Re}\left(\left(\phi_{+}^{-1 / 2} \pm \imath O_{-}^{-1 / 2}\right) b(x ; 0)\right)-\int_{\mathbf{R}} \operatorname{Re} p_{1}\left(o_{\nu}\right) c_{u}^{\prime}(x ; v) d u
$$

where the first term in the right-hand side contributes to the main term of $Q^{1}$. These terms cancel and the rest is equal to

$$
Q^{1}=-2 \int_{\mathbf{R}}\left\langle\nabla \epsilon(o) \operatorname{Re}\left(p_{1}(o) c_{u}^{\prime}(x ; u)\right), \nabla h_{\imath}\right\rangle d v
$$

We estimate the Fourier transform of this integral as in (46). For this we need an estimate for the amplitude. By Step 9, we have $c(x ; u)=(2 u)^{-1}(b(x ; u)-b(x ; 0))$, where $b(x ; u) \doteq 2^{1 / 2} \pi \mathbf{n} a_{b}(x, s)|j \sigma|^{-1 / 2}$, in the second term. Therefore

$$
\max _{u}\left|c_{u}^{\prime}(x ; u)\right| \leq C \max _{u}\left|b_{u u}^{\prime \prime}(x ; u)\right| \leq C\left|a_{\supset}(x, s)\right||j \sigma|^{-1 / 2} \leq C|x-y|^{1 / 2}\left|a_{\supset}(x, s)\right|
$$

since the factor $|j|^{-1 / 2}$ is uniformly bounded and $|\sigma| \leq C / r$ in $\Lambda_{\text {a }}$, otherwise $\nabla h_{,}=0$. Finally

$$
\begin{equation*}
\int_{\Gamma(\tau)}\left|\widehat{Q}^{1}(x ; \omega)\right|^{2} d l \leq \frac{C}{|\omega|^{3}} \tau \exp (-\operatorname{Im} \omega \tau) \tag{49}
\end{equation*}
$$

The inequalities (46), (49) and (48) imply (8).
Step 11. Set $U \doteq \bigcup_{i} U_{i} \subset X \times X$ and take a smooth function $h_{D} \in \mathcal{D}(U)$ that is equal to 1 in $D \times D_{0}$; define $P=h_{D} P_{D}+\left(1-h_{D}\right) P_{\infty}$. An estimate for $P$ follows from (24) and (43). By applying Lemma 5.5 we obtain (7). Next we find

$$
\begin{equation*}
Q \doteq \square P-\delta_{y .0}=\left(h_{D} Q_{D}+R_{D}\right)+\left(1-h_{D}\right) Q_{x} \tag{50}
\end{equation*}
$$

where the term $R_{D}=\left[\square, h_{D}\right] P_{D}+\left[\square .1-h_{D}\right] P_{\infty}=\left[\square . h_{D}\right]\left(P_{D}-P_{x}\right)$ is smooth and supported by $K \cap U$. The operator $P_{D}-P_{\infty}$ can be represented in the form (17). By Lemma 5.5 this implies an estimate for the kernel $\widehat{R}_{D}$ similar to (8). This completes the construction for Lemma 5.1.

## 13. The Helmholtz equation with constant velocity

Step 12. Now we are going to find a kernel $W$ such that the sum $S \doteq \widehat{P}+W$ is a fundamental function. This equation is equivalent to $\square_{\omega} W=-\widehat{Q}$. First we solve the Helmholtz equation with constant velocity

$$
\begin{equation*}
\square_{\omega 0} W_{0}=-\widehat{Q} \tag{51}
\end{equation*}
$$

Lemma 13.1. There exists a solution $W_{0}$ of (51) that is analytic in the set $\{\omega: \operatorname{Im} \omega \geq 0$ and $\omega \neq 0\}$ and satisfies

$$
\begin{equation*}
\int_{t(y, x) \leq \tau}\left|\exp \left(\operatorname{Im} \omega t_{0}(y, x)\right) W_{0}(y, x ; \omega)\right|^{2} d x \leq \frac{1}{|\omega|^{4}}\left(C_{1} \tau+C_{2} \tau^{5}\right), \quad \tau \geq \tau_{0} \tag{52}
\end{equation*}
$$

Proof. We solve this equation by means of the convolution

$$
\begin{equation*}
W_{0}(y, x ; \omega) \stackrel{?}{=}-\int_{X} S_{0}(x-z: w) \widehat{Q}(y \cdot z ; w) d z \tag{53}
\end{equation*}
$$

where $S_{0}$ is the source function for (51). This integral does not, however, converge absolutely. We regularize it, using the special structure (50) and set $\widehat{Q}=\widehat{U}_{D}+\widehat{U}_{\infty}$, where $\widehat{U}_{\infty} \doteq\left(1-h_{D}\right) \widehat{Q}_{\infty}$. The term $\widehat{U}_{D} \doteq h_{D} \widehat{Q}_{D}+\widehat{R}_{D}$ is supported by $D_{0}$, and hence the convolution $W_{D} \doteq S_{0} * \widehat{U}_{D}$ is well defined. To estimate this convolution we use the inequality

$$
\begin{equation*}
\int_{D_{0}}\left|\exp (\operatorname{Im} \omega t(y, x)) \widehat{U}_{D}(y, x: \omega)\right|^{2} d x \leq \frac{C}{\left|\omega^{2}\right|^{3}} \tag{54}
\end{equation*}
$$

which follows from (8) for $\widehat{Q}_{D}$ and $\widehat{R}_{D}$. We have

$$
\begin{aligned}
& \left|\exp \left(\operatorname{Im} \omega t_{0}(y, x)\right) W_{D}(y, x ; \omega)\right|=\left|\int_{D_{0}} \exp \left(\operatorname{Im} \omega t_{0}(y \cdot x)\right) S_{0}(x-z ; \omega) \widehat{U}_{D}(y, z ; \omega) d z\right| \\
& \leq \int_{D_{0}}\left|\exp \left(\operatorname{Im} \omega t_{0}(z, x)\right) S_{0}(x-z: \omega) \exp \left(\operatorname{Im} \omega t_{0}(y, z)\right) \widehat{U}_{D}(y, z ; \omega)\right| d z
\end{aligned}
$$

since $t_{0}(y, x) \leq t_{0}(y, z)+t_{0}(z, x)$. Therefore by Lemma 13.2,

$$
\begin{aligned}
& \int_{t_{0}(y, x) \leq \tau}\left|\exp \left(\operatorname{Im} \omega t_{0}(y, x)\right) W_{D}(y, x ; \omega)\right|^{2} d x \\
& \leq\left(\int_{t_{0}(x . z) \leq \tau+\tau_{0}}\left|\exp \left(\operatorname{Im} \omega t_{0}(x, z)\right) S_{0}(x-z ; \omega)\right| d z\right)^{2} \\
& \times \int_{D_{0}}\left|\exp \left(\operatorname{Im} \omega t_{0}(y . z)\right) \widehat{U}_{D}(y, z ; \omega)\right|^{2} d x
\end{aligned}
$$

where $\tau_{0}$ is the diameter of $D_{0}$ in the metric $g_{0}$. The second factor is bounded by $C|\omega|^{-3}$, due to (54) and the inequality $t_{0}(y, z) \leq t(y, z)$. The first factor is bounded by

$$
\frac{C}{|\omega|}\left(\int_{t_{0}(x, z) \leq \tau+\tau_{0}}|x-z|^{-1 / 2} d z\right)^{2} \leq \frac{C}{|\omega|}\left(\tau+\tau_{0}\right)^{3}
$$

because of (64) and $t_{0}(x, z) \leq \mathbf{n}_{0}|x-z|$. This implies

$$
\int_{t_{0}(y, x) \leq \tau}\left|\exp \left(\operatorname{Im} \omega t_{0}(y, x)\right) W_{D}(y \cdot x: \omega)\right|^{2} d x \leq \frac{C}{|\dot{x}|^{4}}\left(\tau+\tau_{0}+\tau_{1}\right)^{3}
$$

since $t(y, x) \leq t_{0}(y, x)+\tau_{1}$, where $\tau_{1}$ is the $g$-diameter of $D_{0}$.
Lemma 13.2. The inequality $\|a * b\|_{2} \leq\|a\|_{2}\|b\|_{1}$ holds for arbitrary $a \in L_{2}(X)$ and $b \in L_{1}(X)$.

Proof. A proof can be done by the Fourier transform.
Next, we define the convolution of the functions $S_{0}$ and $\hat{U}_{\infty} \doteq\left(1-h_{D}\right) \widehat{Q}_{x}$. For this we will use analytic continuation in the variable $r$. We have

$$
\begin{equation*}
\widehat{Q}_{\infty}=F\left(\square A_{>}\right)+F\left(\square A_{<}\right)+F\left(\square A_{>}\right) . \tag{55}
\end{equation*}
$$

We consider the three terms separately. According to (32) and Lemma 9.1 the term $\left(1-h_{D}\right) F\left(\square A_{>}(y, x, t)\right)$ equals the sum of the functions

$$
\begin{align*}
U_{>}^{n}(y, x ; \omega) & \doteq F\left[(t-\tau(s, r))_{ \pm}^{-1 / 2} e(t-\tau(s, r))(\Delta+\mathbf{m}) a(s, r) h_{>}(\sigma)\right]  \tag{56}\\
& =\exp (\imath \omega \tau(s, r)) q(\omega)(\Delta+\mathbf{m}) a(s, r) h_{>}(\sigma)\left(1-h_{D}\right)
\end{align*}
$$

for $n=1, \ldots, N$, where $q(\omega) \doteq F\left(t_{ \pm}^{-1 / 2} e(t)\right)$ and the right-hand side is a univalent (and smooth) function of $x=x(s, r)$. The sign $\pm 1=\operatorname{sgn} j(s, r)$ is constant for each $n$ and the number $N$ of terms depends on $y, t, s$ and $r$. We wish to regularize the integral

$$
\begin{equation*}
W_{>}^{n}(y, x ; \omega) \stackrel{?}{=}-\int_{X \backslash D_{0}} S_{0}(x-z ; \omega) U_{>}^{n}(y, z: \omega) d z \tag{57}
\end{equation*}
$$

First, note that the integral converges absolutely for $\operatorname{Im} \omega>0$, since $S_{0}(x-z ; \omega)=$ $O(\exp (-\operatorname{Im} \omega|x-z|))$ and $U_{>}^{n}(y, z ; \omega)=O(\exp (-\operatorname{Im} \omega t(y, z)))$ by (56). According to (25) the function $\tau(s, r)$ is linear with respect to $r$, and $a(s, r)$ is of the form $a_{0}(s) r^{-1 / 2}$ by (28) and (29). The function $(\Delta+\mathbf{m}) a$ is also analytic in $r$, as can
be seen from (27). Therefore the right-hand side has analytic continuation to the domain $\{\varrho=r+v \imath: r \geq 0$ and $v \in \mathbf{R}\}$ and this continuation satisfies

$$
\begin{align*}
\left|U_{>}^{n}(y, x ; \omega)\right| & \leq C|q(\omega)| \exp (-\operatorname{Im} \omega \tau(s, \varrho)) \\
& \leq \frac{C}{|\omega|^{3 / 2}} \exp \left(-\operatorname{Im} \omega \tau(s, r)-\operatorname{Re} \omega \mathbf{n}_{0}^{2} v(r)\right) \tag{58}
\end{align*}
$$

since $\tau(s, \varrho)=\tau(s, r)+\mathbf{n}_{0}^{2} v \imath, \tau(s, r)=t(y, x)$. Suppose that Re $\omega \geq 0$. Then the righthand side is bounded in the halfplane $v \geq 0$. First, we make the change of variables $z=z(r, s)$, where $d z=j d r d s, j=j(s, r)$. Take into account that the integrand $S_{0}(x-$ $z ; \omega) U_{>}^{n}(y, z ; \omega) j(s, r)$ admits a holomorphic continuation in $r$.

Shift the domain $X \backslash D_{0}$ to the chain $X_{v}$, defined by the mapping $(s, r) \mapsto z=$ $z(s, r) \doteq-\xi_{0}(s)(r+v(r) \imath)+x_{0}(s)$ from $S^{\prime} \times \mathbf{R}$ to the complex space $X_{\mathbf{C}}=X+\imath X$, where $v(r)$ is a suitable continuous function, provided the function $S_{0}(x-z ; \omega)$ has a bounded holomorphic continuation to a 3-chain $Z_{v}$ in $X_{\mathbf{C}}$ such that $\partial Z_{v}=X \backslash D_{0}$ $X_{v}$. The integrals over $X \backslash D_{0}$ and $X_{v}$ will coincide by the Cauchy theorem. According to Proposition 16.1 this condition is satisfied if $\operatorname{Im} \omega \sqrt{(x-z)^{2}} \geq 0$ for $z \in Z_{v}$, where $w^{2}=w_{1}^{2}+w_{2}^{2}$. This follows from the inequality $\left|\mathbf{n}_{0} \operatorname{Im} \omega-\operatorname{Im} \mathbf{m}\right| \leq C$. Take a positive number $v_{0}$, to be specified later. and a continuous function $v=v(r)$ that vanishes for $0 \leq r \leq r_{0}(x) \doteq 4 \mathbf{n}_{0}^{-1}\left(|x|+\max \left|x_{0}(s)\right|\right)$, and equals $\frac{1}{4}\left(r-r_{0}\right)$ for $r \geq r_{0}$. Set $Z_{v} \doteq\left\{z=-\xi_{0}(s)(r+\tau v(r) z)+x_{0}(s): 0 \leq \tau \leq 1\right\}$ and show that

$$
\begin{equation*}
\operatorname{Re}(x-z)^{2} \geq \operatorname{Im}(x-z)^{2} \geq 0 \quad \text { for } z \in Z_{v} \tag{59}
\end{equation*}
$$

This is obvious for $r \leq r_{0}(x)$, since $v(r)=0$. Otherwise, it follows from the equations

$$
\begin{aligned}
& \operatorname{Re}(x-z)^{2}=\left(\xi_{0} r+x-x_{0}\right)^{2}-\left(\mathbf{n}_{0} \tau v(r)\right)^{2} \\
& \operatorname{Im}(x-z)^{2}=2\langle\operatorname{Re} z-x . \operatorname{Im} z\rangle=2 \tau v(r)\left(\mathbf{n}_{0}^{2} r+\left\langle x-x_{0}, \xi_{0}\right\rangle\right) .
\end{aligned}
$$

This property implies that $Z_{v}$ is contained in the domain $\left\{\omega: \operatorname{Im} \omega \sqrt{(x-z)^{2}} \geq 0\right\}$ for Re $\omega \geq 0$; the function $S_{0}(x-z ; \omega)$ decreases exponentially in $Z_{v}$, as $r \rightarrow \infty$. Then we apply a version of Lemma 13.1 to the integral (57) taken over $X_{v}$. The last factor in (58) is bounded by $\exp (-c \operatorname{Re} \omega r)$ for some $c>0$. We have $\operatorname{Re} \omega>0$ and $r \geq \mathbf{n}_{0}^{-1}|y-z|-r_{1}$ for some $r_{1}$, which ensures convergence of an integral like (57) taken over $X_{v}$.

## 14. Quasianalytic continuation

Next we consider the second term in (55). It is a finite sum of the terms $U_{<}^{n}$ that are similar to (56) with the extra factor $h_{<}(r \sigma), \nabla h_{<}(r \sigma)$ or $\Delta h_{<}(r \sigma)$. This
factor has analytic contimuation for complex $g=r+\imath \imath \in \mathbf{C}$. except for the interval $K(s) \doteq\left\{r: \frac{1}{3} \delta \leq r|\sigma(s)| \leq 3 \mu\right\}$. This set might be unbounded near critical values of the variable $s$. We continue $h$ to the function in the complex plane:

$$
\tilde{h}(\varrho) \doteq \sum_{q=0}^{3} \frac{(\imath v)^{q}}{q!} h^{(q)}(r) . \quad \varrho=r+v \imath
$$

such that $\bar{\partial} \tilde{h}(\varrho)=\frac{1}{2} v^{3} h^{(4)}(r)=O\left(v^{3}\right)$ and also

$$
\begin{equation*}
\partial^{l} \bar{\partial}^{k+1} \tilde{h}(\varrho)=O\left(v^{3-l-k}\right) . \quad l+k \leq 2 . \tag{60}
\end{equation*}
$$

Now we regularize the integral

$$
W_{<}^{n}(x ; \omega) \stackrel{?}{=}-\int_{X \backslash D_{0}} S_{0}(x-z: w) U_{<}^{n}(y, z ; w) d z
$$

by means of the Cauchy-Green theorem as follows:

$$
\begin{aligned}
W_{<}^{n}(x ; \omega) \doteq & -\int_{X_{v}} S_{0}(x-z: \omega) U_{<}^{n}(y . z: \omega) j d r d s \\
& -\int_{Z_{v}} S_{0}(x-z: \omega) \bar{\partial}_{\underline{\varrho}} U_{<}^{n}(y . z: \omega) j d \underline{\varrho} d \underline{\varrho} d s
\end{aligned}
$$

where $Z_{v}=\{z=z(r+v u, s): 0 \leq v \leq v(r)\} \subset X_{\mathbf{C}}$. The first integral is similar to that of the function $W_{>}^{n}$. The inequality $|x-\operatorname{Re} z| \leq \operatorname{Re} \sqrt{(x-z)^{2}}$ holds for $z \in Z_{v}$, by (59). On the other hand, Re $\omega \operatorname{Im} \sqrt{(x-z)^{2}} \geq 0$. This yields

$$
\operatorname{Im} \omega|x-y| \leq \operatorname{Im} \omega|x-\operatorname{Re} z|+\operatorname{Im} \omega|y-\operatorname{Re} z| \leq \operatorname{Im} \omega \sqrt{(x-z)^{2}}+\operatorname{Im} \omega|y-\operatorname{Re} z|
$$

By this inequality

$$
\begin{aligned}
\exp \left(\operatorname{Im} \omega \mathbf{n}_{0}|x-y|\right) \mid \int_{Z_{v}} & S_{0}(x-z: \dot{v}) \bar{\partial}_{\underline{\varrho}} U_{<}^{n}\left(y \cdot z ; \omega^{\prime}\right) j d \underline{\varrho} d \underline{\varrho} d s \mid \\
\leq & \sup _{z \in Z_{v}}\left|\exp \left(\mathbf{n}_{0} \operatorname{Im} \omega \sqrt{(x-z)^{2}}\right) S_{0}(x-z ; \omega)\right| \\
& \times \int_{Z_{v}} \exp \left(\operatorname{Im} \omega \mathbf{n}_{0}|y-\operatorname{Re} z|\right)\left|\bar{\partial}_{\underline{\varrho}} U_{<}^{n}(y, z ; \omega) j\right| d v d r d s \\
\leq & \frac{C}{|\omega|^{1 / 2}} \int_{Z_{v}} \exp \left(\operatorname{Im} \omega \mathbf{n}_{0}|y-\operatorname{Re} z|\right)\left|\bar{\partial}_{\varrho} U_{<}^{n}(y, z ; \omega) j\right| d v d r d s
\end{aligned}
$$

The supremum is finite by Proposition 16.1. Applying an inequality like (58) to $U_{<}^{n}$ and taking into account that $\operatorname{Re} \omega \mathbf{n}_{0}^{2} v \geq 0$, we get

$$
\begin{align*}
& |\omega|^{3 / 2} \int_{Z_{v}} \exp \left(\operatorname{Im} \omega \mathbf{n}_{0}|y-\operatorname{Re} z|\right)\left|\bar{\partial}_{\underline{g}} U_{<}^{n}(y \cdot z: \omega) j\right| d v d r d s \\
& \leq C \int_{Z_{v}}\left(\left|j \bar{\partial}_{\underline{g}} \tilde{h}\right|+\left|j \bar{\partial}_{\underline{g}} \nabla \tilde{h}^{\prime}\right|+\left|j \bar{\partial}_{\underline{g}} \Delta \tilde{h}\right|\right) d v d r d s \tag{61}
\end{align*}
$$

where $\tilde{h}=\tilde{h}(\sigma \varrho)$. To estimate the first term in the right-hand side, we use the relations $\bar{\partial}_{\varrho} \tilde{h}(\sigma \varrho)=O\left(\sigma^{4}|v|^{3}\right), j=\varrho \sigma+O(1)$. Then

$$
\begin{aligned}
\int_{Z_{v}}\left|j \bar{\partial}_{\varrho} \tilde{h}\right| d v d r d s & \leq C_{1} \int_{S^{1}}|\sigma|^{5} \int_{K} \int_{0}^{v(r)}|v|^{4}|j| d v d r d s \\
& \leq C_{2} \int_{S^{1}}|\sigma|^{5} \int_{K(s)}\left(\left(\log r / r_{0}\right)^{6}+C\right) d r d s \leq C_{3}
\end{aligned}
$$

since $\sigma^{3} \log \left(r / r_{0}\right)$ is uniformly bounded, $r \in K(s)$. and $s \in S^{1}$. For the third term in (61) we use (27) and (60) which gives $\left|j \bar{\partial}_{\underline{g}} \Delta \tilde{h}(\sigma \underline{Q})\right| \leq C|\sigma|^{4} \underline{\varrho}^{2} v$. This estimate is sufficient for the finiteness of the integral of $\left|j \bar{\partial}_{\underline{g}} \Delta \tilde{h}\right| d v d r d s$ over $Z_{v}$. The second and third terms in (61) are estimated similarly.

Now we have only to consider the third term in (55). We have

$$
\square_{\omega} \hat{A}_{b}(x ; \omega)=\int_{\mathbf{R}} \exp \left(\imath \omega \varphi_{j}(x, s)\right) \operatorname{Re} q_{1}(\omega)(\Delta+\mathbf{m})\left(a,(x, s) h_{j}(r, s)\right) d s
$$

where $q_{1}(\omega)=F\left(p_{1}(t)\right)$ and the omitted terms contain derivatives of $h_{\nu}$. The phase function $\varphi_{b}(x, s)=t_{0}(s)+\mathbf{n}_{0}^{2} r$ for $r \doteq \mathbf{n}_{0}^{-2}\left\langle\xi_{0}(s) . x_{0}(s)-x\right\rangle$ admits analytic continuation with respect to $r$ and $\operatorname{Im} \varphi_{j}(x, s)=\mathbf{n}_{0}^{2} r$. This function is nonnegative if $v \geq 0$. According to the choice (38) the function $a_{;}(x, s)$ does not depend on $r$ at all. Therefore only the $\bar{\partial}_{r}$-derivative of the factors (47) appear in $\bar{\partial}_{r} \square_{i}, \hat{A}_{3}$. We estimate the corresponding integrals as in the previous case. In this way we define the regularized convolution $W_{b}=S_{0} * \square_{\omega} \hat{A}_{3}$. Taking the sum $W_{0}=\sum W_{>}^{n}+\sum W_{<}^{n}+W_{b}$, we get a solution of (51) that admits an analytic continuation at the domain $\{\omega: \operatorname{Im} \omega \geq 0$ and $\operatorname{Re} \omega>0\}$, which satisfies (52). To get an analytic continuation to the domain Re $\omega<0$, we replace $v$ by $-v$ in the above arguments. This completes the proof of Lemma 13.1.

## 15. The completion of the construction

Step 13. We have $\square_{\omega}\left(\left(1-h_{D}\right) W_{0}\right)+\widehat{Q}=-L$. where $h_{D} \in \mathcal{D}(X)$ is as in Step 11. By (51) the kernel $L=h_{D} \widehat{Q}-\left\langle\nabla h_{D}, \nabla W_{0}\right\rangle-\Delta h_{D} W_{0}$ is supported in $D \times D_{0}$. We
take into account (52) and estimate the derivatives of $W_{0}$ by the standard energy integral method

$$
\int_{D_{0}}\left|\exp \left(\operatorname{Im} \omega t_{0}(y \cdot x)\right) \nabla W_{0}(y \cdot x: \omega)\right|^{2} d x \leq \frac{C}{|\omega|^{2}}
$$

This and (8) imply

$$
\begin{equation*}
\int_{X}\left|\exp \left(\operatorname{Im} \omega t_{0}(y, x)\right) L(y \cdot x ; \omega)\right|^{2} d x \leq \frac{C_{L}}{|\dot{\omega}|^{2}} \tag{62}
\end{equation*}
$$

Now we look for a kernel $R$ such that $S=\widehat{P}+W_{0}+R$ is a fundamental solution. It is now equivalent to $\square_{\omega} R=L$. The kernel $T \doteq \widehat{P}+W_{0}$ satisfies $\square_{\omega} T=\delta_{y}-L$. By (7),

$$
\int_{t(y, x) \leq \tau}\left|\exp \left(\operatorname{Im} \omega t_{0}(y, x)\right) T(y, x ; \omega)\right|^{2} d x \leq \frac{C_{T}}{|\dot{\omega}|}\left(\tau^{5}+\tau\right)
$$

The sum $S=T+R$ is a fundamental function, if $R$ is a solution of the integral equation

$$
R(y, x ; \omega)-\int_{X} L(y, z ; \omega) R(z, x ; \omega) d z=\int_{X} L(y, z ; \omega) T(z, x ; \omega) d z
$$

Write the solution by means of the Neumann series $R \doteq R^{(1)}+R^{(2)}+\ldots$ where

$$
R^{(k)}(y, x ; \omega) \doteq \int_{X} \ldots \int_{X} L\left(y, z_{k} ; \omega\right) \ldots L\left(z_{3} . z_{2} ; \omega\right) L\left(z_{2}, z_{1} ; \omega\right) T\left(z_{1}, x ; \omega\right) d z_{k} \ldots d z_{1}
$$

For an arbitrary $y \in D$, we have by (62).

$$
\begin{align*}
\int_{X} \mid \exp \left(\operatorname{Im} \omega t_{0}(y, x)\right) & \left.R^{(k)}(y, x ; \omega)\right|^{2} d x \\
\leq & \int_{X} \mid \exp \left(\operatorname{Im} \omega t_{0}(y, x)\right) \\
& \times\left.\int_{X} \ldots \int_{X} L\left(y, z_{k} ; \omega\right) \ldots L\left(z_{2}, z_{1}: \omega\right) T\left(z_{1}, x: \omega\right) d z_{k} \ldots d z_{1}\right|^{2} d x \\
\leq & \int_{X}\left|\exp \left(\operatorname{Im} \omega t_{0}(y, z)\right) L(y, z: \omega)\right|^{2} d z  \tag{63}\\
& \times\left(\int_{D} \int_{D}\left|\exp \left(\operatorname{Im} \omega t_{0}(z, x)\right) L(z \cdot x ; \omega)\right|^{2} d z d x\right)^{k-1} \\
& \times \int_{X}\left|\exp \left(\operatorname{Im} \omega t_{0}(z, x)\right) T(z, x ; \omega)\right|^{2} d z d x
\end{align*}
$$

$$
\begin{aligned}
& \leq \frac{C_{T}}{|\omega|}\left(V(D) \int_{X}\left|\exp \left(\operatorname{Im} \omega t_{0}(y, z)\right) L(y, z: \omega)\right|^{2} d z\right)^{k} \\
& \leq \frac{C_{T}\left(C_{L} V(D)\right)^{k}}{|\omega|^{2 k+1}}
\end{aligned}
$$

where $V(D)$ stands for the volume of $V$, since $t_{0}(y, x) \leq t_{0}\left(y, z_{k}\right)+\ldots+t_{0}\left(z_{1}, x\right)$. The series converges if $|\omega|^{2}>C_{L} V(D)$. and we obtain the estimate for $|\dot{\omega}| \geq \omega_{0} \doteq$ $\left(2 C_{L} V(D)\right)^{1 / 2}$ :

$$
\int_{X}\left|\exp \left(\operatorname{Im} \omega t_{0}(y, x)\right) R\left(y . x ; \omega^{\prime}\right)\right|^{2} d x \leq \frac{C_{T}}{V(D)|\omega|^{3}} .
$$

The inequality (3) follows from (7) and (63). This completes the construction of the source function $S$ that satisfies (3) for all $y \in D$. The compact $D$ can be magnified, consequently the above arguments are valid for arbitrary compact subsets $D$ of the plane. The sequence of the corresponding source functions gives the global function in $X \times X$ because of the uniqueness theorem. To complete the proof of Theorem 3.1 we only need to check that the inequality (3) is satisfied with some constants $C_{1}$ and $C_{2}$ that do not depend on $y \in X$. Starting from a point $y \in X \backslash D$, in the first step we use the source function $S_{0}$ for the homogeneous space. It satisfies (3) with $C_{2}=0$ according to Proposition 16.1 of the next section. Moreover. the arc integrals of the density $\left|\exp \left(\operatorname{Im} \omega t_{0}(y, x)\right) S_{0}(y, x ; \omega)\right|^{2} d l$ over $\Gamma(\tau)$ are uniformly bounded. Indeed, if $y$ is at the distance $r$ from $D$, only the portion $O\left(r^{-1}\right)$ of the energy is scattered on the inhomogeneity of the medium. This, together with the arguments of Steps $2-13$ shows that the estimate (3) is uniform for $y \in X$. This completes the proof of Theorem 3.1.

## 16. The source function in homogeneous medium

Denote by $\log \zeta$ the univalent branch that is defined for $|\arg \zeta|<\pi$ and is positive for $\zeta>1$. Set $\sqrt{\zeta}=\exp \left(\frac{1}{2} \log \zeta\right)$.

Proposition 16.1. If $\mathbf{p}$ is constant and $\operatorname{Im} \mathbf{p} \geq 0$, then the source function $S_{0}$ of the operator $\Delta+\mathbf{p}^{2}$ has the form

$$
S_{0}(x ; \omega)=\frac{1}{4 \pi} J_{0}\left(\sqrt{(\mathbf{p} x)^{2}}\right) \log (\mathbf{p} x)^{2}+A\left((\mathbf{p} x)^{2}\right)
$$

where $J_{0}$ is the Bessel function and $A$ is an entire analytic function of order $\frac{1}{2}$. Moreover, it satisfies the inequality

$$
\begin{equation*}
\left|S_{0}(x ; \omega)\right| \leq 8 \pi\left|\sqrt{(\mathbf{p} x)^{2}}\right|^{-1 / 2} \exp \left(-\operatorname{Im} \mathbf{p} \sqrt{x^{2}}\right) \tag{64}
\end{equation*}
$$

in the domain $\left|\arg (\mathbf{p} x)^{2}\right|<\pi$, of the complex 2-space $X_{\mathbf{C}}$.

Proof. For real $\mathbf{p}$ the distribution $E_{0}(x . t)=-(2 \pi)^{-1} \theta(t-\mathbf{p}|x|)\left(t^{2}-\mathbf{p}^{2}|x|^{2}\right)^{-1 / 2}$ defined in $X \times \mathbf{R}$ satisfies the equation $\left(\Delta-\mathbf{p}^{2} \partial_{t}^{2}\right) E_{0}=\delta_{0}$. Set $\varrho=\sqrt{(\mathbf{p} x)^{2}}$ and calculate

$$
\begin{aligned}
&-2 \pi S_{0}(x)=\int_{t \geq r} \exp (t \imath)\left(t^{2}-\varrho^{2}\right)^{-1 / 2} d t \stackrel{t=\underline{\varrho} s}{\underline{=}} \int_{s \geq 1} \exp (\varrho s \imath)\left(s^{2}-1\right)^{-1 / 2} d s \\
& s=\stackrel{1+\sigma t}{=}-\imath \exp (\underline{\varrho}) \int_{0}^{\infty} \exp (-\underline{\varrho} \sigma)\left(-\sigma^{2}+2 \sigma \imath\right)^{-1 / 2} d \sigma .
\end{aligned}
$$

The right-hand side admits an analytic continuation to the halfplane Re $\varrho>0$. The second factor is estimated by means of the Cauchy theorem as follows:

$$
\begin{aligned}
\left|\int_{0}^{\infty} \exp (-\varrho \sigma)\left(-\sigma^{2}+2 \sigma \imath\right)^{-1 / 2} d \sigma\right| & =\left|\int_{0}^{\infty} \exp (-\underline{\varrho} \sigma) \sigma^{-1 / 2}(2 \imath-\sigma)^{-1 / 2} d \sigma\right| \\
& \stackrel{\sigma=\overline{\bar{\varrho}} \tau}{=} \underline{\varrho}\left|\int_{0}^{\infty} \exp \left(-|\varrho|^{2} \tau\right)(\bar{\varrho} \tau)^{-1 / 2}(2 \imath-\bar{\varrho} \tau)^{-1 / 2} d \tau\right| \\
& \leq \max _{\tau}|2 \imath-\bar{\varrho} \tau|^{-1 / 2}|\varrho|^{1 / 2}\left|\int_{0}^{\infty} \frac{\exp \left(-|\varrho|^{2} \tau\right)}{\tau^{1 / 2}} d \tau\right| \\
& =\left(\frac{1}{2} \pi\right)^{1 / 2}|\varrho|^{-1 / 2} .
\end{aligned}
$$

Remark. We apply this proposition to the coefficient $\mathbf{p}=\sqrt{\omega^{2} \mathbf{n}_{0}^{2}+\mathbf{m}_{0}}$ for positive $\mathbf{n}_{0}$. Therefore if $\operatorname{Im} \omega \geq 0$ and $|\omega| \geq\left|\mathbf{m}_{0}\right|^{1 / 2} \mathbf{n}_{0}^{-1}$. the coefficient $\mathbf{p}$ is also in the upper halfplane.

## 17. The quantitative stationary phase method

Introduce a norm in $\mathbf{R}^{n}$ and denote by $\|a\|$ the norm of an operator $a$ in $\mathbf{R}^{n}$.
Proposition 17.1. Let $x_{0} \in \mathbf{R}^{n}$ be a critical point of a smooth phase function $\varphi$ and $U$ be a starlike neighbourhood of $x_{0}$ such that

$$
\begin{equation*}
\left\|\left\{\varphi_{i j}^{\prime \prime}(x)-\psi_{i j}^{\prime \prime}\left(x_{0}\right)\right\}\right\|\left\|\left\{\gamma_{i j}^{\prime \prime}\left(x_{0}\right)\right\}^{-1}\right\|<\frac{1}{2} \tag{65}
\end{equation*}
$$

for $x \in U$. Then there exists a smooth coordinate system $y=y(x)$ in $U$ such that $y\left(x_{0}\right)=0$ and $\beta(y(x))=\varphi(x)$, where $\beta(y)=\frac{1}{2} \psi_{i j}^{\prime \prime}\left(x_{0}\right) y^{i} y^{j}$.

Proof. Assume that $x_{0}=0$ and write $\boldsymbol{\varphi}(x)-3(x)=\sum_{i . j=1}^{n} a_{i j}(x) x^{i} x^{j}$, where

$$
\begin{equation*}
a_{i j}(x)=\int_{0}^{1} \int_{0}^{1}\left[\hat{\gamma}_{i j}^{\prime \prime}(s t x)-\psi_{i j}^{\prime \prime}(0)\right] d s d t \tag{66}
\end{equation*}
$$

Look for a solution of form $y(x)=x+z x$, where $z=z(x)$ is an unknown matrix function vanishing at the origin, and $x$ is written as a column. The equation for $z$ is $\beta(x, z x)+\beta(z x, x)+\beta(z x, z x)=x^{\prime} a x$. where $x^{\prime}$ means the row of coordinate functions and $a=\left\{a_{i j}\right\}$. Therefore it is sufficient to solve the matrix equation $z+b^{-1} z^{\prime} b+b^{-1} z^{\prime} b z=c \doteq b^{-1} a$, where $b=\left\{\frac{1}{2} \psi_{i j}^{\prime \prime}\left(x_{0}\right)\right\}$. Since the matrices $b$ and $a$ are symmetric, we have $b^{-1}\left(c^{\prime}\right)^{k} b=c^{k}$ for any natural $k$. Write $z$ as a power series in $c$ such that $z(0)=0$. The system takes the form $2 z+z^{2}=c$ and we find $z=z(c)=(1+c)^{1 / 2}-1$. The series converges if $\|c(x)\|<1$. We have by (66) and (65) that

$$
\|c(x)\| \leq\left\|b^{-1}\right\|\|a(x)\| \leq 2\left\|\left\{\varphi_{i j}^{\prime \prime}\left(x_{0}\right)\right\}^{-1}\right\|\left\|\left\{\gamma_{i j}^{\prime \prime}(x)-\gamma_{i j}^{\prime \prime}\left(x_{0}\right)\right\}\right\|<1 .
$$

Proposition 17.2. Let $\beta(x)=\frac{1}{2} b_{i j} x^{i} x^{j}$ be a nondegenerated form of signature $s$ and $f$ be a continuous function such that $f . g \in L_{1}\left(\mathbf{R}^{n}\right)$. The integral

$$
I(\omega)=\int_{\mathbf{R}^{n}} \exp (2 \pi \imath \dot{\omega}(x)) f(x) d x . \quad w^{\prime}>0
$$

can be written in the form

$$
I(\omega)=\frac{\exp \frac{1}{4} \imath \pi s}{\omega^{n / 2} \sqrt{\operatorname{det} b+0 \imath}} f(0)+\frac{1}{\pi \omega \imath} \int_{\mathbf{R}^{n}} \exp (2 \pi \imath \dot{\omega} 3(x)) g(x) d x
$$

where

$$
\begin{equation*}
g \doteq-\sum_{j=1}^{n} \alpha^{j}(D) f_{j}, \quad \sum_{j=1}^{n} x^{j} f_{j}(x)=f(x)-f(0) . \quad a^{j}=a^{j k} \partial_{k} . \quad\left\{a^{j k}\right\}=\left\{b_{j k}\right\}^{-1} \tag{67}
\end{equation*}
$$

Proof. A proof can be done by the regularization $3 \mapsto 3+\imath \varepsilon x^{2}$ and partial integration.

An estimate can be done for the function $g$ :
Proposition 17.3. If $U \subset \mathbf{R}^{n}$ is an arbitrary starlike neighbourhood of the origin, then there exists a function $g$, that satisfies (67) and $\max _{U}|g| \leq \max _{U}|\alpha(D) f|$. where $\alpha(D)=a^{j k} \partial_{j} \partial_{k}$.

Proof. A proof follows from the equations

$$
\begin{aligned}
f_{j}(x) & =\int_{0}^{1} \partial_{j} f(t x) d t \\
g(x) & =-\sum_{j=1}^{n} \int_{0}^{1} \alpha^{j}(D) \partial_{j} f(t x) t d t=-\int_{0}^{1} \alpha(D) f(t x) t d t
\end{aligned}
$$

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