# Asymptotics without logarithmic terms for crack problems 

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Dedicated to Prof. W.L. WendLand for his $65^{\text {th }}$ anniversary.


#### Abstract

We consider boundary value problems for elliptic systems in a domain complementary to a smooth surface $\mathscr{M}$ with boundary $\mathscr{E}$. The same boundary conditions are prescribed on both sides of the surface $\mathscr{M}$. The most important model behind this investigation is the crack problem in three-dimensional linear elasticity (either isotropic or anisotropic): there the boundary conditions are Neumann, i.e. tractions are prescribed on both faces of the crack surface $\mathscr{M}$. We prove that the singular functions appearing in the expansion of the solution along the crack edge $\mathscr{E}$ all have the form $r^{\frac{1}{2}+k} \psi(\theta)$ in local polar coordinates $(r, \theta)$ : the logarithmic shadow terms predicted by the general theory do not appear. Moreover, we obtain that, for a smooth right hand side, the jump of the displacement across the crack surface is the product of $r^{\frac{1}{2}}$ with a smooth vector function on $\mathscr{M}$.

We present two different, but complementing, approaches leading to these results, and providing distinct generalizations. The first one is based on a Wiener-Hopf factorization of the pseudodifferential symbol on the surface $\mathscr{M}$ obtained after reduction of the boundary value problem. The condition on the symbol which yields the absence of logarithmic terms into the solution of the boundary pseudodifferential equation is a variant of the transmission condition. The asymptotics of the solution in the full space is then deduced by a representation formula from the asymptotics of the solution on $\mathscr{M}$. The second approach concerns directly the boundary value problem and is based on a closer look at the Mellin symbol at each point of the crack edge $\mathscr{E}$. The Mellin symbol is proved to act between special subspaces of angular functions and the absence of logarithmic terms is the consequence of a series of compatibility conditions, valid for any Agmon-Douglis-Nirenberg system.


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## Part A. Scope and principal results

## A. 1 The Crack domain and the boundary value problem

Let $\mathscr{M}$ be a bounded $\mathscr{C}^{\infty}$ orientable surface of codimension 1 in $\mathbb{R}^{n+1}$. We assume that $\mathscr{M}$ is a manifold with $\mathscr{C}^{\infty}$ boundary $\mathscr{E}$. We consider boundary value problems set in the domain

$$
\Omega:=\mathbb{R}^{n+1} \backslash \overline{\mathscr{M}} .
$$

For the equations of linear elasticity (Lamé or, more generally, anisotropic material law), the solutions of such boundary value problems yield the stresses in the domain $\Omega$ around $\mathscr{M}$ which represents a crack with front $\mathscr{E}$. For the equations of electromagnetism (Helmholtz or Maxwell), the solutions represent the diffracted field around the screen $\mathscr{M}$.

We are going to set our problem and describe our results in a framework including such problems, which is also covered by the hypotheses of our two methods.

We denote by $\boldsymbol{x}=\left(x_{1}, \ldots, x_{n+1}\right)$ cartesian coordinates in $\mathbb{R}^{n+1}$ and by $\partial_{\boldsymbol{x}}^{\alpha}$ the partial derivative $\partial_{1}^{\alpha_{1}} \ldots \partial_{n+1}^{\alpha_{n+1}}$. Let $b$ be a homogeneous integrodifferential form of degree 1 with constant coefficients acting on $N$ component vectors $\boldsymbol{u}, \boldsymbol{v} \in \mathrm{H}^{1}(\Omega)^{N}$

$$
b(\boldsymbol{u}, \boldsymbol{v})=\sum_{j=1}^{N} \sum_{k=1}^{N} \sum_{|\alpha|,|\beta|=1} \int_{\Omega} a_{j k}^{\alpha \beta} \partial_{\boldsymbol{x}}^{\alpha} u_{j} \partial_{\boldsymbol{x}}^{\beta} \bar{v}_{k} \mathrm{~d} \boldsymbol{x}
$$

Here $\boldsymbol{u}=\left(u_{1}, \ldots, u_{N}\right), \boldsymbol{v}=\left(v_{1}, \ldots, v_{N}\right)$ and the coefficients $a_{j, k}^{\alpha, \beta}$ are constant. We assume that the form $b$ is coercive on $\mathrm{H}^{1}(\Omega)^{N}$, i.e. that for some constants $c, C>0$ there holds

$$
\begin{equation*}
\forall \boldsymbol{u} \in \mathrm{H}^{1}(\Omega)^{N}, \quad \operatorname{Re} b(\boldsymbol{u}, \boldsymbol{u})+C\|\boldsymbol{u}\|_{\mathrm{L}^{2}(\Omega)}^{2} \geq c\|\boldsymbol{u}\|_{\mathrm{H}^{1}(\Omega)}^{2} . \tag{A1}
\end{equation*}
$$

Moreover, we suppose that $b$ is symmetric on $\mathrm{H}^{1}(\Omega)^{N}$ :

$$
\begin{equation*}
\forall \boldsymbol{u}, \boldsymbol{v} \in \mathrm{H}^{1}(\Omega)^{N}, \quad b(\boldsymbol{u}, \boldsymbol{v})=\overline{b(\boldsymbol{v}, \boldsymbol{u})} . \tag{A2}
\end{equation*}
$$

The partial differential operator associated with the bilinear form $b$ is

$$
L=\left(L_{j k}\right)_{j, k} \quad \text { with } \quad L_{j k}=-\sum_{|\alpha|,|\beta|=1} \partial_{x}^{\beta} a_{j k}^{\alpha \beta} \partial_{x}^{\alpha} .
$$

Hypotheses $\left(\mathfrak{H}_{\mathrm{A} 1}\right)$ and $\left(\mathfrak{H}_{\mathrm{A} 2}\right)$ are satisfied for the Laplace equation $(N=1)$, for the equations of general elasticity, including the anisotropic case, ( $N$ is equal to the dimension of the space) and for equations of thermoelasticity and electroelasticity ( $N$ is the dimension of the space plus 1 ).

Since $\mathscr{M}$ is orientable, we can define a smooth unit normal vector field $\boldsymbol{n}$ on $\mathscr{M}$, which is unique if we choose the direction of the normal at some fixed point. After fixing the field $\boldsymbol{n}$ we can fix the traces $\gamma_{ \pm}$, taking $\gamma_{+}$opposite to the direction of $\boldsymbol{n}$ (i.e. from "above" if we consider $\boldsymbol{n}$ as pointing upward) and taking $\gamma_{-}$in the direction of $\boldsymbol{n}$ (i.e. from "below").

The Neumann operator $T$ associated with $b$ and the normal field $\boldsymbol{n}$ is defined as

$$
T=\left(T_{j k}\right)_{j, k} \quad \text { with } \quad T_{j k}=\sum_{|\alpha||,|\beta|=1} \boldsymbol{n}^{\beta} a_{j k}^{\alpha \beta} \partial_{\boldsymbol{x}}^{\alpha}, \quad \boldsymbol{n}^{\beta}=n_{1}^{\beta_{1}} \ldots n_{n+1}^{\beta_{n+1}} .
$$

Let $B$ denote either the identity (which will be associated with the Dirichlet operator) or the Neumann operator $T$ on $\mathscr{M}$. We consider solutions $\boldsymbol{u} \in \mathrm{H}^{1}(\Omega)^{N}$ of the problem

$$
\left\{\begin{array}{rll}
L \boldsymbol{u} & =\boldsymbol{f} & \text { in } \Omega  \tag{A.1.1}\\
\gamma_{ \pm} B \boldsymbol{u} & =0 & \\
\text { on } \mathscr{M},
\end{array}\right.
$$

with, possibly, conditions at infinity (note that we may relax the condition $\boldsymbol{u} \in \mathrm{H}^{1}(\Omega)^{N}$ into $\boldsymbol{u} \in \mathrm{H}^{1}\left(\Omega \cap \mathscr{B}_{R}\right)^{N}$ for any $R>0$, with $\mathscr{B}_{R}$ the ball of center 0 and radius $R$ ). We assume that $f$ is a $\mathscr{C}^{\infty}$ vector function on $\mathbb{R}^{n+1}$, with compact support.

## A. 2 State of the art and motivations

Due to the presence of the edge $\mathscr{E}$, the domain is highly non-smooth and this yields strong singularities for the solutions of problem (A.1.1) along this edge. The general structure of these singularities is known, and addressed by many works, see [ChkDu2, CoDa1, Da1, DuWe1, Gr1, KozMaRo1, MaPl1, MaRo1]. The generic form of these singularities is

$$
c\left(\mathscr{X}^{\prime}\right) \sum_{q=0}^{Q} r^{\lambda\left(x^{\prime}\right)} \log ^{q} r \boldsymbol{\psi}_{q}\left(\mathscr{X}^{\prime}, \theta\right)
$$

where $\mathscr{X}^{\prime}$ represents coordinates in $\mathscr{E}$, and $(r, \theta)$ polar coordinates in the planes normal to $\mathscr{E}$, centered on $\mathscr{E}$.

The structure of $\boldsymbol{u}$ in a neighborhood of the boundary $\mathscr{E}$ of $\mathscr{M}$ is very important in applications. For example, in elasticity, it provides an essential tool for the investigation of crack propagation in the quasi-static case. The propagation criterion is based on the stress intensity factors (the coefficients $c\left(x^{\prime}\right)$ of the leading terms in asymptotics) and on the "polarization operator" (which involves the second terms in asymptotics), see [Na1].

Thus, it is important to know that the asymptotic expansion contains neither oscillatory terms (i.e. non real exponents $\lambda$ ) nor logarithmic terms (ie. $\log ^{q} r$ with $q \geq 1$ ). Concerning oscillations, it is known that the solution of the crack problems never oscillates provided the crack is inside an homogeneous material, even if the material is anisotropic, see [DuWe1]. Concerning logarithms, although absence of logarithms in the leading terms was known long
ago for isotropic materials [Gr2, NaPl1], the same was not proved for further terms, where logarithms could appear as a shadow singularities.

The main scope of the present investigation is to establish that the structure of the solution $\boldsymbol{u}$ of the general problem (A.1.1) is simpler than the general theory would predict. The main result can be summarized in one sentence:
"The edge asymptotics of $\boldsymbol{u}$ does not contain any logarithmic term $\log r$."
Still in the framework of elasticity, this was observed in the case of a curved crack in the isotropic elastic plane $\mathbb{R}^{2}$ for the second term in the asymptotics in [WeSt1, Theorem 2.4] and in the case of a half plane crack $\mathbb{R}_{+}^{2}$ in the anisotropic elastic space $\mathbb{R}^{3}$ in [DuWe1, Theorem 4.3]; For curved cracks the conjecture was first formulated by S.A. NAZAROV.

Moreover, it has been shown [CoDa4, DuNa1], that even in the very general framework of Agmon-Douglis-Nirenberg systems with the same boundary conditions on both sides of the crack $\mathscr{M}$, the principal part of the asymptotics contains only powers of $r$ with halfinteger exponents (ie. $\lambda=\frac{1}{2}+k, k \in \mathbb{N}_{0}$ ), and without any $\log r$ term, see also [Koz1] for scalar operators of order $2 m$ with Dirichlet condition. In this work, we prove that, in fact, this simple structure extends to the complete asymptotics.

The result that the whole asymptotics does not contain $\log r$ terms is by no way obvious, and is not an easy consequence of the simple structure of its principal part. Indeed, because the exponents $\frac{1}{2}+k$ of the whole asymptotics are all translated from each other by integers, we should expect $\log r$ terms, due to the interaction between the non-principal terms in the operator and the principal singularities (see, for example, [KozMaRo1, Remark 10.5.1], where this interaction is explained).

## A. 3 REDUCTION TO THE CRACK SURFACE AND REPRESENTATION FORMULAS

One of the essential features of our crack-type boundary value problem (A.1.1) is that all information on the singular behavior of $\boldsymbol{u}$ is contained in an $N$-component vector function $\phi$, defined on the crack surface by the jump of $\boldsymbol{u}$ across $\mathscr{M}$

$$
\boldsymbol{\phi}=[C \boldsymbol{u}]:=\gamma_{+} C \boldsymbol{u}-\gamma_{-} C \boldsymbol{u}
$$

where $C$ denotes the complementing trace of $B$, i.e. the Dirichlet trace if $B$ is Neumann and $C=T$ if $B$ is Dirichlet. Of course, the asymptotics of $\boldsymbol{u}$ will yield the asymptotics of $\phi$. But even more important is that $\phi$ can be directly obtained as the solution of a pseudodifferential equation on $\mathscr{M}$ of the form

$$
\begin{equation*}
\mathbf{a}\left(\mathscr{x}, D_{\mathscr{X}}\right) \boldsymbol{\phi}(\mathscr{x})=\boldsymbol{g}(x), \quad x \in \mathscr{M}, \tag{A.3.1}
\end{equation*}
$$

and analyzed in this respect. The relation between the boundary value problem (A.1.1) and the pseudodifferential equation (A.3.1) will be fully explained in § B.2. Let us only mention that in the case of the Dirichlet problem $\mathbf{a}=V$, where $V:=\gamma_{+} \mathscr{V}=\gamma_{-} \mathscr{V}$ is the trace of
the single layer potential $\mathscr{V}$ associated with the operator $L$, and in the case of the Neumann problem, $\mathbf{a}=W$, where $W:=\gamma_{-} T \mathscr{D}=\gamma_{+} T \mathscr{D}$ is the Neumann trace of the double layer potential $\mathscr{D}$. Then $\boldsymbol{u}$ can be reconstructed by the representation formula

$$
\begin{equation*}
\forall \boldsymbol{x} \in \Omega, \quad \boldsymbol{u}(\boldsymbol{x})=N \boldsymbol{f}(\boldsymbol{x})+\mathscr{D}[\boldsymbol{u}](\boldsymbol{x})-\mathscr{V}[T \boldsymbol{u}](\boldsymbol{x}) \tag{A.3.2}
\end{equation*}
$$

where $[\boldsymbol{u}]:=\gamma_{+} \boldsymbol{u}-\gamma_{-} \boldsymbol{u}$ and $[T \boldsymbol{u}]:=\gamma_{+} T \boldsymbol{u}-\gamma_{-} T \boldsymbol{u}$ denote the jumps of the functions $\boldsymbol{u}(\boldsymbol{x})$ and $T \boldsymbol{u}(\boldsymbol{x})$ across the surface $\mathscr{M}$ and $N$ denotes the Newton (volume) potential. Thus, the asymptotics of $\boldsymbol{u}$ depends only on $\phi$ because the volume potential part $N \boldsymbol{f}$ is smooth and $[B \boldsymbol{u}]=0$ on $\mathscr{M}$.

Note that the coerciveness hypothesis $\left(\mathfrak{H}_{\mathrm{A} 1}\right)$ ensures the Fredholm property of both problems (A.1.1) and (A.3.1) in appropriate spaces.

Thus, two different approaches are available to us: either first study the solution $\phi$ of equation (A.3.1), then derive the asymptotics of $\boldsymbol{u}$, or first study the solution $\boldsymbol{u}$ of problem (A.1.1), then derive the asymptotics of $\phi=[C \boldsymbol{u}]$.

First approach. The first approach is exposed in Part B: we develop the potential operator technique based on the Wiener-Hopf factorization, according to the three main following steps:
St. 1 The boundary value problem (A.1.1) is reduced to a pseudodifferential equation of type (A.3.1) on the crack surface $\mathscr{M}$ by invoking the representation of solutions (A.3.2), see $\S$ B.2.
St. 2 Asymptotics of solutions $\phi$ of the pseudodifferential equation on the crack surface are found using the Wiener-Hopf factorization, see § B. 3 - § B.7.
St. 3 By inserting the surface asymptotics into the representation formula (A.3.2), the full spatial asymptotic expansion of $\boldsymbol{u}$ is derived, see $\S$ B.8.
G. Eskin was the first who applied the Wiener-Hopf factorization to investigations of asymptotics, see [Es1]. The method received contributions by several authors [Be1, ChkDu1, ChkDu2, ChkDu3, CoSt1, DuWe1, DuNa1]. In particular in [ChkDu1] necessary and sufficient conditions for the absence of logarithms in the principal part of the asymptotics were found. Here we obtain criteria for the absence of logarithms in the whole asymptotics.

SECOND APPROACH. The second approach is exposed in Part C: it relies on the classical Mellin transform, $c f[\mathrm{Ko1}]$, and more recent representation formulas for the angular part of singular functions, of [CoDa2]. The main steps are:
St. 1 By separation of variables and Mellin transform in $r$, the problem is transformed into systems of ordinary differential equations in the angular variable $\theta$ with the parameters $\mathscr{X}^{\prime}$ and $\lambda$, the dual variable of $r$, see $\S$ C.1.
St. 2 The solutions of these systems are represented by contour integrals around the unit circle with the Cayley symbols of the principal part of the operator, see §C.4.
St. 3 By the Cayley representation formulae, the condition of absence of logarithm is reduced to compatibility conditions between traces of a series of right hand sides in the Mellin calculus, see $\S$ C.5.

## A. 4 Results

In order to state our results, let us introduce local coordinates in a neighborhood of the edge $\mathscr{E}$ which is the crack front.

Definition A.4.1 (i) Let $\mathscr{X}^{\prime}=\left(\mathscr{X}_{1}, \ldots, \mathscr{X}_{n-1}\right)$ denote local coordinates in $\mathscr{E}$.
(ii) For $\mathscr{X}^{\prime} \in \mathscr{E}$, let $\Pi_{\mathscr{X}}$, denote the normal plane to $\mathscr{E}$ containing $\mathscr{X}^{\prime}$. We take polar coordinates $(r, \theta)$ in $\Pi_{\mathscr{X}^{\prime}}$ such that $r=0$ is the intersection $\Pi_{\mathscr{X}^{\prime}} \cap \mathscr{E}, \theta=-\pi$ is $\Pi_{\mathscr{X}^{\prime}} \cap \mathscr{M}$ from below and $\theta=\pi$ is $\Pi_{\mathscr{X}^{\prime}} \cap \mathscr{M}$ from above.
(iii) We set $x_{n}=r \cos \theta$ and $x_{n+1}=r \sin \theta$. The $n$ coordinates $\left(\mathscr{x}^{\prime}, x_{n}\right)$ are local coordinates in $\mathscr{M}$ and the $n+1$ coordinates $\mathscr{x}:=\left(\mathscr{X}^{\prime}, \mathscr{X}_{n}, \mathscr{X}_{n+1}\right)$ are local coordinates in $\Omega$ in a neighborhood of $\mathscr{E}$.
(iv) The local cylindrical coordinates are $\left(\mathscr{X}^{\prime}, r, \theta\right)$ and we shall use $\left(\mathscr{X}^{\prime}, r, 0\right)=\mathscr{X} \in \mathscr{M}$ and $\left(\mathscr{X}^{\prime}, 0,0\right)=\mathscr{X}^{\prime} \in \mathscr{E}$.
(v) The dual variables of $\underline{\mathscr{x}}=\left(\mathscr{X}^{\prime}, \mathscr{X}_{n}, \mathscr{X}_{n+1}\right)$ are denoted by $\underline{\xi}=\left(\xi^{\prime}, \xi_{n}, \xi_{n+1}\right)$.
(vi) We denote by $\kappa: \underline{x} \mapsto \boldsymbol{x}$ the generic map of an atlas on $\mathscr{M}$, and by $\mathscr{J}_{\kappa}(\underline{x}):=$ $\left[D \kappa(\underline{x})^{\top}\right]^{-1}$, the inverse of its Jacobian matrix.

From the combination of general edge asymptotics [MaPl1, MaRo1, NaPl1], [Da1, CoDa1], and of the particular structure of the principal part for crack problems [DuWe1, CoDa4], we may derive that there holds the following general statement, see $\S$ B. 2 and C.2.

## Proposition A.4. 2

(i) Any solution $\boldsymbol{u}$ of the boundary value problem (A.1.1) with a smooth right hand side $f$ has the following asymptotic expansion as $r \rightarrow 0$ : For any integer $K \geq 0$

$$
\begin{array}{r}
\boldsymbol{u}=\sum_{j=1}^{N} c_{j}^{0}\left(\mathscr{X}^{\prime}\right) r^{\frac{1}{2}} \boldsymbol{\psi}_{j}^{0}\left(\mathscr{X}^{\prime}, \theta\right)+\sum_{k=1}^{K} \sum_{q=0}^{q(k)} \sum_{j=1}^{j(k)} c_{j}^{k, q}\left(\mathscr{X}^{\prime}\right) r^{\frac{1}{2}+k} \log ^{q} r \boldsymbol{\psi}_{j}^{k, q}\left(\mathscr{X}^{\prime}, \theta\right)  \tag{A.4.1}\\
+\boldsymbol{u}_{\mathrm{reg}, K}+\boldsymbol{u}_{\mathrm{rem}, K}
\end{array}
$$

The coefficients $c_{j}^{0}, c_{j}^{k, q}$ are $\mathscr{C}^{\infty}(\mathscr{E})$ functions depending on $\boldsymbol{f}$. The regular part $\boldsymbol{u}_{\mathrm{reg}, K}$ is a linear combination of terms of the form $c\left(\mathscr{x}^{\prime}\right) p\left(x_{n}, x_{n+1}\right)$, with polynomial $p$ and $\mathscr{C}^{\infty}(\mathscr{E})$ coefficient c. The remainder $\boldsymbol{u}_{\mathrm{rem}, K}$ satisfies $\partial^{\beta} \boldsymbol{u}_{\mathrm{rem}, K}=\mathscr{o}\left(r^{K-|\beta|+1 / 2}\right)$ as $r \rightarrow 0$ for any multi-index $\beta \in \mathbb{N}_{0}^{n+1}$. The $\boldsymbol{\psi}_{j}^{0}$ and $\boldsymbol{\psi}_{j}^{k, q}$ are $N$-component vector functions in $\mathscr{C}^{\infty}([-\pi, \pi] \times \mathscr{E})$ and depend only on the domain $\Omega$ and the operators $(L, B)$.
(ii) Any solution $\phi=[C \boldsymbol{u}]$ of the pseudodifferential equation (A.3.1) with a smooth right hand side $\boldsymbol{g}$ has the following asymptotic expansion as $r \rightarrow 0$ : For any integer $K$

$$
\begin{equation*}
\phi=r^{\frac{\nu}{2}} \boldsymbol{d}^{0}\left(x^{\prime}\right)+\sum_{k=1}^{K} \sum_{q=0}^{q(k)} r^{\frac{\nu}{2}+k} \log ^{q} r \boldsymbol{d}^{k, q}\left(\mathscr{X}^{\prime}\right)+\boldsymbol{\phi}_{\mathrm{rem}, K} \tag{A.4.2}
\end{equation*}
$$

Here $\nu$ is the order of the pseudodifferential operator a. The $\boldsymbol{d}^{0}$ and $\boldsymbol{d}^{k, q}$ are $N$ component vector functions in $\mathscr{C}^{\infty}(\mathscr{E})$. The remainder $\phi_{\mathrm{rem}, K}$ satisfies $\partial^{\beta} \phi_{\mathrm{rem}, K}=$ $\sigma\left(r^{K-|\beta|+\nu / 2}\right)$ as $r \rightarrow 0$ for any multi-index $\beta \in \mathbb{N}_{0}^{n+1}$.

Our main result in this paper is that there are no logarithmic terms at all in expansions (A.4.1) and (A.4.2):

## Theorem A.4.3

(i) Any solution $\boldsymbol{u}$ of the boundary value problem (A.1.1) with smooth right hand side $\boldsymbol{f}$ has the following asymptotic expansion as $r \rightarrow 0$ : For any integer $K \geq 0$

$$
\begin{array}{r}
\boldsymbol{u}=\sum_{j=1}^{N} c_{j}^{0}\left(\mathscr{X}^{\prime}\right) r^{\frac{1}{2}} \boldsymbol{\psi}_{j}^{0}\left(\mathscr{X}^{\prime}, \theta\right)+\sum_{k=1}^{K} \sum_{j=1}^{j(k)} c_{j}^{k}\left(x^{\prime}\right) r^{\frac{1}{2}+k} \boldsymbol{\psi}_{j}^{k}\left(\mathscr{X}^{\prime}, \theta\right)  \tag{A.4.3}\\
+\boldsymbol{u}_{\mathrm{reg}, K}+\boldsymbol{u}_{\mathrm{rem}, K} .
\end{array}
$$

The scalar coefficients $c_{j}^{k}$ belong to $\mathscr{C}^{\infty}(\mathscr{E})$ and depend on $\boldsymbol{f}$, while the $N$-component vector functions $\boldsymbol{\psi}_{j}^{k}$ depend only on the domain $\Omega$ and the operators $(L, B)$.
(ii) Any solution $\boldsymbol{\phi}=[C \boldsymbol{u}]$ of the pseudodifferential equation (A.3.1) with smooth right hand side $\boldsymbol{g}$ has the following asymptotic expansion as $r \rightarrow 0$ : For any integer $K \geq 0$

$$
\begin{equation*}
\boldsymbol{\phi}=r^{\frac{\nu}{2}} \boldsymbol{d}^{0}\left(x^{\prime}\right)+\sum_{k=1}^{K} r^{\frac{\nu}{2}+k} \boldsymbol{d}^{k}\left(x^{\prime}\right)+\boldsymbol{\phi}_{\mathrm{rem}, K} . \tag{A.4.4}
\end{equation*}
$$

The coefficients $\boldsymbol{d}^{0}$ and $\boldsymbol{d}^{k, q}$ are $N$-component vector functions in $\mathscr{C}^{\infty}(\mathscr{E})$.
This result is proved in Parts B and C, in more general frameworks: In Part B, for a general class of pseudo-differential equations with classical symbols satisfying a "continuity property" which is a sort of variant of the transmission condition. In Part C, it is proved for a large class of Agmon-Douglis-Nirenberg systems with covering boundary conditions. Moreover both approaches allow precise representation formulas for the "angular" vector functions $\psi_{j}^{k}\left(\mathscr{X}^{\prime}, \theta\right)$ as linear combinations of simple trigonometric functions, see $\S$ B. 8 and C.7.

Because of the relation $\phi=[C \boldsymbol{u}]$ between $\boldsymbol{u}$ and $\phi$, it is quite simple to link the first terms in expansions (A.4.3) and (A.4.4).

In the Neumann case, $C=\mathrm{Id}$ and:

$$
\boldsymbol{d}^{0}\left(\mathscr{X}^{\prime}\right)=\sum_{j=1}^{N} c_{j}^{0}\left(\mathscr{X}^{\prime}\right)\left[\boldsymbol{\psi}_{j}^{0}\left(\mathscr{X}^{\prime}, \theta\right)\right]_{\pi},
$$

where $[\boldsymbol{\psi}(\theta)]_{\pi}$ denotes the jump $\boldsymbol{\psi}(\pi)-\boldsymbol{\psi}(-\pi)$. In the Dirichlet case, $C=T$, and let $r^{-1} T_{0}\left(\mathscr{X}^{\prime}, \theta ; r \partial_{r}, \partial_{\theta}\right)+T_{1}\left(\mathscr{X}^{\prime}\right) \partial_{\mathscr{X}^{\prime}}$ be the expression of $T$ in cylindrical coordinates. Then there holds

$$
\boldsymbol{d}^{0}\left(\mathscr{X}^{\prime}\right)=\sum_{j=1}^{N} c_{j}^{0}\left(\mathscr{X}^{\prime}\right)\left[T_{0}\left(\mathscr{X}^{\prime}, \theta ; \frac{1}{2}, \partial_{\theta}\right) \boldsymbol{\psi}_{j}^{0}\left(\mathscr{X}^{\prime}, \theta\right)\right]_{\pi} .
$$

Defining $s_{j}^{0}\left(\mathscr{X}^{\prime}\right) \in \mathscr{C}^{\infty}(\mathscr{E}) \otimes \mathbb{C}^{N}$ by

$$
\begin{cases}\boldsymbol{s}_{j}^{0}\left(\mathscr{X}^{\prime}\right)=\left[\boldsymbol{\psi}_{j}^{0}\left(\mathscr{X}^{\prime}, \theta\right)\right]_{\pi} & \text { if } C=\mathrm{Id} \\ \boldsymbol{s}_{j}^{0}\left(\mathscr{X}^{\prime}\right)=\left[T_{0}\left(\theta, \mathscr{X}^{\prime} ; \frac{1}{2}, \partial_{\theta}\right) \boldsymbol{\psi}_{j}^{0}\left(\mathscr{X}^{\prime}, \theta\right)\right]_{\pi} & \text { if } C=T\end{cases}
$$

we get the common relation

$$
\begin{equation*}
\boldsymbol{d}^{0}\left(x^{\prime}\right)=\sum_{j=1}^{N} c_{j}^{0}\left(x^{\prime}\right) s_{j}^{0}\left(x^{\prime}\right) \tag{A.4.5}
\end{equation*}
$$

The vectors $s_{j}^{0}\left(\mathscr{X}^{\prime}\right), j=1, \ldots, N$ are independent of the right hand side and form a basis of $\mathbb{C}^{N}$ for each fixed $\mathscr{X}^{\prime}$. We will address in a forthcoming paper formulae and numerical methods for computing the scalar coefficients $c_{j}^{0}\left(x^{\prime}\right)$.

Conversely, as a consequence of the representation formula (A.3.2), we obtain the inverse relation between the coefficients involved in (A.4.5): all coefficients $c_{j}^{0}\left(x^{\prime}\right)$ are defined as a composition of some matrices with $\boldsymbol{d}^{0}\left(\mathscr{X}^{\prime}\right)$, see [ChkDu2].

## A. 5 MODULAR REPRESENTATION

The asymptotics (A.4.3) and (A.4.4) give the possibility of representing $\boldsymbol{u}$ and $\phi$ as finite linear combination of non-smooth functions with smooth coefficients: As a straightforward consequence of (A.4.4), we obtain the following factorization of the density $\phi$ :

Corollary A.5.1 Any solution $\phi$ of the boundary pseudodifferential equation (A.3.1) with a smooth right hand side $\boldsymbol{g}$ satisfies

$$
\begin{equation*}
r^{-\frac{\nu}{2}} \phi \in \mathscr{C}^{\infty}(\overline{\mathscr{M}})^{N} \tag{A.5.1}
\end{equation*}
$$

As a further consequence of the expansion (A.4.3), we can prove that a simple splitting of $\boldsymbol{u}$ holds in local cylindrical coordinates. For this, we first introduce $\mathscr{U}$, a closed tubular neighborhood of the edge $\mathscr{E}$ where the local cartesian coordinates are well defined. We may take $\mathscr{U}$ as a set of the form

$$
\mathscr{U}=\left\{\left(\mathscr{X}^{\prime}, \mathscr{X}_{n}, \mathscr{X}_{n+1}\right) ; r \leq r_{0}, \mathscr{X}^{\prime} \in \mathscr{E}\right\} .
$$

Then we denote by $\breve{\mathscr{U}}$ its expression in local cylindrical coordinates

$$
\stackrel{\mathscr{U}}{ }=\left\{\left(\mathscr{X}^{\prime}, r, \theta\right) ; 0 \leq r \leq r_{0}, \theta \in[-\pi, \pi], \mathscr{X}^{\prime} \in \mathscr{E}\right\} .
$$

Note that we clearly distinguish the two faces $\theta=-\pi$ and $\theta=\pi$ of $\breve{\mathscr{U}}$.
Corollary A.5.2 Let $\boldsymbol{u}$ be any solution of the problem (A.1.1) with a smooth right hand side $\boldsymbol{f}$ and denote by $\breve{\boldsymbol{u}}$ its expression in local cylindrical coordinates: $\boldsymbol{u}\left(\mathscr{X}^{\prime}, \mathscr{X}_{n}, \mathscr{X}_{n+1}\right)=$ $\breve{\boldsymbol{u}}\left(\mathscr{X}^{\prime}, r, \theta\right)$. Then $\breve{\boldsymbol{u}}$ admits a splitting in two parts

$$
\begin{equation*}
\breve{\boldsymbol{u}}\left(\mathscr{X}^{\prime}, r, \theta\right)=\breve{\boldsymbol{u}}_{0}\left(\mathscr{X}^{\prime}, r, \theta\right)+r^{\frac{1}{2}} \breve{\boldsymbol{u}}_{1}\left(\mathscr{X}^{\prime}, r, \theta\right), \tag{A.5.2}
\end{equation*}
$$

where $\breve{\boldsymbol{u}}_{0}$ and $\breve{\boldsymbol{u}}_{1}$ are $\mathscr{C}^{\infty}(\breve{\mathscr{U}})$ in the variables $r, \theta$ and $\mathscr{X}^{\prime}$.

Now, we may write (A.5.2) in local cartesian coordinates and obtain

$$
\begin{equation*}
\boldsymbol{u}\left(\mathscr{X}^{\prime}, \mathscr{X}_{n}, \mathscr{X}_{n+1}\right)=\boldsymbol{u}_{0}\left(\mathscr{X}^{\prime}, \mathscr{X}_{n}, \mathscr{X}_{n+1}\right)+r^{\frac{1}{2}} \boldsymbol{u}_{1}\left(\mathscr{X}^{\prime}, \mathscr{X}_{n}, \mathscr{X}_{n+1}\right) . \tag{A.5.3}
\end{equation*}
$$

The part $u_{0}$ in (A.5.3) is in fact $\mathscr{C}^{\infty}(\overline{\mathscr{U}})$ in the coordinates $\left(\mathscr{X}^{\prime}, \mathscr{X}_{n}, \mathscr{X}_{n+1}\right)$. Now we may wonder if $\boldsymbol{u}_{1}$ is also a $\mathscr{C}^{\infty}(\overline{\mathscr{U}})$ function. This is not true. For example, for the Laplace operator with Dirichlet boundary conditions we have $u_{1}=c_{1} \sin \frac{1}{2} \theta+c_{2} r \sin \frac{3}{2} \theta+\cdots$ Replacing the factor $r^{\frac{1}{2}}$ by another function does not help. We need to split $r^{1 / 2} u_{1}$ into new parts. Again, when $L=\Delta$ and $n=2$, we simply have

$$
r^{\frac{1}{2}} u_{1}=c_{1}\left(\zeta^{\frac{1}{2}}-\bar{\zeta}^{\frac{1}{2}}\right)+c_{2}\left(\zeta^{\frac{3}{2}}-\bar{\zeta}^{\frac{3}{2}}\right)+c_{3}\left(\zeta^{\frac{5}{2}}-\bar{\zeta}^{\frac{5}{2}}\right) \cdots
$$

with $\zeta=r e^{i \theta}$. Therefore

$$
r^{\frac{1}{2}} u_{1}=\zeta^{\frac{1}{2}}\left(c_{1}+c_{2} \zeta+c_{3} \zeta^{2}+\cdots\right)+\bar{\zeta}^{\frac{1}{2}}\left(c_{1}+c_{2} \bar{\zeta}+c_{3} \bar{\zeta}^{2}+\cdots\right)
$$

which means that $r^{1 / 2} u_{1}$ can be written as $\zeta^{\frac{1}{2}} u_{1}^{\prime}+\bar{\zeta}^{\frac{1}{2}} u_{2}^{\prime}$ with $\mathscr{C}^{\infty}(\overline{\mathscr{U}})$ functions $u_{1}^{\prime}$ and $u_{2}^{\prime}$. This result extends to the wider class of problems satisfying hypotheses ( $\mathfrak{H}_{\mathrm{A} 1}$ ) and $\left(\mathfrak{H}_{\mathrm{A} 2}\right)$, provided a condition on the symbol of the interior operator $L$ : the symbol $\boldsymbol{\xi} \mapsto L(\boldsymbol{\xi})$ of $L$ is defined so that $L=L\left(D_{x}\right)$, where $D_{x}=i \partial_{x}$. We require that this symbol satisfies,
$\left(\mathfrak{H}_{\text {A3 }}\right) \quad \forall \mathscr{X}^{\prime} \in \mathscr{E}, \quad$ the roots $\tau \in \mathbb{C}$ of $\quad \operatorname{det} L\left(\mathscr{J}_{\kappa}\left(\mathscr{X}^{\prime}\right)(0,1, \tau)\right)=0 \quad$ are simple,
where we recall, $c f$ Definition A.4.1, that $x^{\prime}$ stands for $x=\left(x^{\prime}, 0,0\right)$ and $(0,1, \tau)$ is the value of the dual variable $\underline{\xi}=\left(\xi^{\prime}, \xi_{n}, \xi_{n+1}\right)$. Note that $\bar{L}\left(\mathscr{J}_{k}\left(\mathscr{X}^{\prime}\right) \underline{\xi}\right)$ is the principal part of the symbol of the operator $L$ written in local variables $(\underline{x} ; \underline{\xi})$.

Theorem A.5.3 If hypotheses $\left(\mathfrak{H}_{\mathrm{A} 1}\right)-\left(\mathfrak{H}_{\mathrm{A} 3}\right)$ are satisfied, then there exist $2 N$ scalar singular functions $\sigma_{\ell}=r^{1 / 2} \varphi_{\ell}\left(\mathscr{X}^{\prime}, \theta\right)$ for $\ell=1, \ldots, 2 N$, with $\varphi_{\ell} \in \mathscr{C}^{\infty}(\mathscr{E} \times[-\pi, \pi])$ such that any solution $\boldsymbol{u}$ of the problem (A.1.1) with smooth right hand side $f$ can be split as follows

$$
\begin{equation*}
\boldsymbol{u}=\boldsymbol{u}_{0}+\sigma_{1} \boldsymbol{u}_{1}^{\prime}+\cdots+\sigma_{2 N} \boldsymbol{u}_{2 N}^{\prime} \tag{A.5.4}
\end{equation*}
$$

where $\boldsymbol{u}_{0}, \boldsymbol{u}_{1}^{\prime}, \ldots, \boldsymbol{u}_{2 N}^{\prime}$ are $\mathscr{C}^{\infty}(\overline{\mathscr{U}})$-smooth vector functions in local cartesian variables.

## Part B. The Wiener-Hopf approach

In this part we investigate the asymptotics of solutions of a class of Pseudo-Differential Equations ( $\Psi \mathrm{DE}$ ) on the manifold $\mathscr{M}$; we also study how these asymptotics are transformed by representation formulas and how they give back asymptotics for our class of Boundary Value Problems (BVP).

In $\S$ B.1, we fix notations for more or less classical Sobolev and Bessel potential spaces, including anisotropic Bessel potential spaces.

In $\S$ B.2, we recall how the boundary value problem (A.1.1) with the Dirichlet or Neumann boundary conditions can be reduced to the $\Psi \mathrm{DE}$ (A.3.1) on the manifold $\mathscr{M}$. The feedback is governed by the representation formulas which reconstruct the solution of the BVP in $\Omega$ from the solution of the $\Psi$ DE on $\mathscr{M}$.

In $\S$ B.3, we introduce a large class of classical $\Psi D E$ on $\mathscr{M}$ and recall from [ChkDu1, Es1] the general form of asymptotics of the solutions $\phi$ of such equations near the boundary $\mathscr{E}$ of $\mathscr{M}$.

In § B.4, we concentrate our attention on a sub-class of classical $\Psi D E$ where the full symbol satisfies a special continuity condition, denoted $\left(\mathfrak{H}_{\mathrm{B} 4}\right)$, with respect to the conormal variable and state the main result of Part B: the asymptotics of the solutions do not contain any logarithmic term (see Theorem B.4.1). We prove that the $\Psi D E$ (A.3.1) obtained from the BVP (A.1.1) belong to our sub-class of $\Psi \mathrm{DE}$.

In $\S$ B.5, before proving the main theorem in its full general framework, we investigate the simpler situation of scalar $\Psi$ DO in dimension 1 . We find a necessary and sufficient condition, denoted $\left(\mathfrak{H}_{\mathrm{B} 5}\right)$ for the absence of logarithms from the whole asymptotics: the continuity condition ( $\mathfrak{H}_{\mathrm{B} 4}$ ) we exhibit in the general situation of dimension $n$ for systems appears as a particular case of $\left(\mathfrak{H}_{\mathrm{B} 5}\right)$.

In § B.6, we give useful auxiliary propositions relating to $\Psi D O$ in one variable acting on functions of $n$ variables and in $\S$ B.7, we prove the main Theorem B.4.1.

In $\S$ B.8, relying on results from [ChkDu2], we give, as a consequence of the simple structure of the solutions $\phi$ of $\Psi D E$, the form of vector functions $\boldsymbol{u}$ defined in $\Omega$ by a certain type of representation formula acting on $\phi$. We prove that the representation formulae (A.3.2) belong to this type. As a result we have the statement of Theorem A.4.3.

## B. 1 Sobolev and Bessel potential spaces

## B.1.A Standard spaces

We first recall the definition of the Fourier transform and Sobolev spaces.
Let $\mathscr{S}\left(\mathbb{R}^{n+1}\right)$ denote the Schwartz space of all rapidly decreasing functions and $\mathscr{S}^{\prime}\left(\mathbb{R}^{n+1}\right)$ the dual space of tempered distributions. For $\varphi \in \mathscr{S}^{\prime}\left(\mathbb{R}^{n+1}\right)$ let

$$
\mathscr{F} \varphi(\boldsymbol{\xi})=\mathscr{F}_{\boldsymbol{y} \rightarrow \boldsymbol{\xi}} \varphi(\boldsymbol{\xi}):=\int_{\mathbb{R}^{n+1}} e^{i \boldsymbol{\xi} \cdot \boldsymbol{y}} \varphi(\boldsymbol{y}) \mathrm{d} \boldsymbol{y}, \quad \boldsymbol{\xi} \in \mathbb{R}^{n+1}
$$

denote its Fourier transform in $\mathbb{R}^{n+1}$. The inverse Fourier transform $\mathscr{F}_{\xi \rightarrow y}^{-1}$ in $\mathbb{R}^{n+1}$ is defined as

$$
\mathscr{F}_{\boldsymbol{\xi} \rightarrow \boldsymbol{y}}^{-1} \psi(x):=\frac{1}{(2 \pi)^{n+1}} \int_{\mathbb{R}^{n+1}} e^{-i \boldsymbol{y} \cdot \boldsymbol{\xi}} \psi(\xi) \mathrm{d} \boldsymbol{\xi}
$$

We denote by $\mathscr{F}_{x \rightarrow \xi}$ and $\mathscr{F}_{\xi \rightarrow x}^{-1}$ the Fourier and inverse Fourier transforms in $\mathbb{R}^{n}$.
The Sobolev space $H^{s}\left(\mathbb{R}^{n+1}\right)$ is defined as the subspace of $\mathscr{S}^{\prime}\left(\mathbb{R}^{n+1}\right)$ endowed with the norm

$$
\|\varphi\|_{\mathrm{H}^{s}\left(\mathbb{R}^{n+1}\right)}^{2}:=\int_{\mathbb{R}^{n+1}}\left(1+|\boldsymbol{\xi}|^{2}\right)^{s}\left|\mathscr{F}_{y \rightarrow \boldsymbol{\xi}} \varphi(\boldsymbol{\xi})\right|^{2} \mathrm{~d} \boldsymbol{\xi}
$$

For an integer $s=m \in \mathbb{N}_{0}$ an equivalent norm on the space $\mathrm{H}^{m}\left(\mathbb{R}^{n+1}\right)$ is

$$
\left(\sum_{|\alpha| \leq m} \int_{\mathbb{R}^{n+1}}\left|\partial_{\boldsymbol{y}}^{\boldsymbol{\alpha}} \varphi(\boldsymbol{y})\right|^{2} \mathrm{~d} \boldsymbol{y}\right)^{1 / 2}
$$

For a domain $\Omega \subset \mathbb{R}^{n+1}$ with a smooth boundary ( $\Omega$ can be, for example, one of the half-spaces $\mathbb{R}_{ \pm}^{n+1}:=\mathbb{R}^{n} \times \mathbb{R}^{ \pm}$), two families of spaces can be defined:
(i) The subspace $\widetilde{\mathrm{H}}^{s}(\Omega) \subset \mathrm{H}^{s}\left(\mathbb{R}^{n+1}\right)$ of the distributions $\varphi$ which are supported inside $\bar{\Omega}$. The extension by 0 outside $\bar{\Omega}$ of such a distribution yields an element in $\mathrm{H}^{s}\left(\mathbb{R}^{n+1}\right)$.
(ii) The quotient space $\mathrm{H}^{s}(\Omega):=\mathrm{H}^{s}\left(\mathbb{R}^{n+1}\right) / \widetilde{\mathrm{H}}^{s}\left(\Omega^{c}\right)$, where $\Omega^{c}:=\mathbb{R}^{n+1} \backslash \bar{\Omega}$ is the complementary domain. The space $\mathrm{H}^{s}(\Omega)$ can also be interpreted as the space of restrictions $p_{\Omega} \varphi$ of functions $\varphi \in \mathrm{H}^{s}\left(\mathbb{R}^{n+1}\right)$. The space is endowed with the factornorm, i.e. the minimal norm of all possible extensions to $\mathbb{R}^{n+1}$.
By $\widetilde{\mathrm{H}}^{s}(\Omega)^{N}, \mathrm{H}^{s}(\Omega)^{N}$, we will denote the spaces of $N$-vector functions.
For a surface $\mathscr{M} \subset \mathbb{R}^{n+1}$ of codimension 1 , with a smooth boundary $\partial \mathscr{M}$, the spaces $\mathrm{H}^{s}(\mathscr{M})$ and $\widetilde{\mathrm{H}}^{s}(\mathscr{M})$ are defined in a standard way, involving some fixed finite covering $\left\{U_{j}\right\}_{j=1}^{J}$ of $\mathscr{M}$, appropriate diffeomorphisms $\varkappa_{j}: U_{j} \rightarrow V_{j} \subset \mathbb{R}_{+}^{n}$ and partition of a unity subordinate to the fixed covering, see, e.g. [Es1, Hr 1$]$.

## B.1.B ANISOTROPIC WEIGHTED SPACES

Besides the above classical spaces, we need a 3-parameter class of anisotropic Sobolev spaces with weight. The weight appears as integer powers of one particular coordinate. We first define these spaces on $\mathbb{R}^{n}$, then on $\mathbb{R}_{+}^{n}$, finally on $\mathscr{M}$.

Let $\mu, s \in \mathbb{R}$ and $\kappa \in \mathbb{N}_{0}$. We denote by $\mathrm{H}^{(\mu, s), \kappa}\left(\mathbb{R}^{n}\right)$ the Hilbert space of distributions $u$ with finite norm

$$
\|u\|_{\mathrm{H}^{(\mu, s), \kappa}\left(\mathbb{R}^{n}\right)}^{2}:=\sum_{k=0}^{\kappa}\left\|\left\langle D^{\prime}\right\rangle^{\mu}\langle D\rangle^{s+k} x_{n}^{k} u\right\|_{L_{2}\left(\mathbb{R}^{n}\right)}^{2} \simeq \sum_{k=0}^{\kappa}\left\|\left\langle\xi^{\prime}\right\rangle^{\mu}\langle\xi\rangle^{s+k} \mathscr{F}\left[D_{n}^{k} u\right]\right\|_{L_{2}\left(\mathbb{R}^{n}\right)}^{2}
$$

where $x=\left(x^{\prime}, x_{n}\right)$ are cartesian coordinates in $\mathbb{R}^{n}, D_{n}:=i \partial_{n}, \xi=\left(\xi^{\prime}, \xi_{n}\right)$ are the corresponding dual variables,

$$
\langle\xi\rangle:=\left(1+|\xi|^{2}\right)^{\frac{1}{2}},
$$

and where

$$
\left\langle D^{\prime}\right\rangle^{\mu}:=\mathscr{F}_{\xi^{\prime} \rightarrow x^{\prime}}^{-1}\left\langle\xi^{\prime}\right\rangle^{\mu} \mathscr{F}_{y^{\prime} \rightarrow \xi^{\prime}}, \quad\langle D\rangle^{\nu}:=\mathscr{F}_{\xi \rightarrow x}^{-1}\langle\xi\rangle^{\nu} \mathscr{F}_{y \rightarrow \xi} .
$$

are the Bessel potential operators. For integer $\mu, s \in \mathbb{N}_{0}$ we have the equivalent norm

$$
\sum_{k=0}^{\kappa} \sum_{\substack{\alpha^{\prime} \in \mathbb{N}^{n-1} \\\left|\alpha^{\prime}\right| \leq \mu}} \sum_{\substack{\beta \in \mathbb{N}^{n} \\|\beta| \leq s+k}}\left\|\partial_{x^{\prime}}^{\alpha^{\prime}} \partial_{x}^{\beta}\left[x_{n}^{k} u\right]\right\|_{L_{2}\left(\mathbb{R}^{n}\right)}
$$

We define the Frechet spaces

$$
\mathrm{H}^{(\infty, s), \kappa}\left(\mathbb{R}^{n}\right):=\bigcap_{\mu \in \mathbb{N}} \mathrm{H}^{(\mu, s), \kappa}\left(\mathbb{R}^{n}\right) \quad \text { and } \quad \mathrm{H}^{(\infty, s), \infty}\left(\mathbb{R}^{n}\right):=\bigcap_{\kappa \in \mathbb{N}} \mathrm{H}^{(\infty, s), \kappa}\left(\mathbb{R}^{n}\right)
$$

The functions in these spaces are $\mathrm{H}^{s}$ globally on $\mathbb{R}^{n}$ and $\mathscr{C}^{\infty}$ in $\mathbb{R}^{n} \backslash\left(\mathbb{R}^{n-1} \times\{0\}\right)$.
On the half-space $\mathbb{R}_{+}^{n}=\mathbb{R}^{n-1} \times \mathbb{R}_{+}$, we define $\mathrm{H}^{(\mu, s), \kappa}\left(\mathbb{R}_{+}^{n}\right)$ as the space of restrictions to $\mathbb{R}_{+}^{n}$ of distributions in $\mathrm{H}^{(\mu, s), \kappa}\left(\mathbb{R}^{n}\right)$. The space $\widetilde{\mathrm{H}}^{(\mu, s), \kappa}\left(\mathbb{R}_{+}^{n}\right)$ denotes the subspace of $\mathrm{H}^{(\mu, s), \kappa}\left(\mathbb{R}^{n}\right)$ of distributions with support in $\overline{\mathbb{R}}_{+}^{n}$.

The spaces $\mathrm{H}^{(\mu, s), \kappa}(\mathscr{M})$ and $\widetilde{\mathrm{H}}^{(\mu, s), \kappa}(\mathscr{M})$ for a smooth compact manifold $\mathscr{M}$ with a smooth boundary $\partial \mathscr{M}$ are defined in a standard way, involving some fixed finite covering of $\mathscr{M}$, appropriate diffeomorphisms and partition of a unity subordinate to the covering, so that the particular coordinate $x_{n}$ corresponds to the distance to $\partial \mathscr{M}$ in $\mathscr{M}$, see [ChkDu1, § 1.1].

## B. 2 REDUCTION TO THE BOUNDARY

In this section, we explain in more detail the way from BVP (A.1.1) to $\Psi D E$ (A.3.1) and back.

We start from the first Green formula for all $\boldsymbol{u} \in \mathrm{H}^{2}(\Omega)^{N}$ and $\boldsymbol{v} \in \mathrm{H}^{1}(\Omega)^{N}$ :

$$
\begin{equation*}
b(\boldsymbol{u}, \boldsymbol{v})=-\int_{\Omega} L \boldsymbol{u} \cdot \overline{\boldsymbol{v}} \mathrm{~d} \boldsymbol{y}+\int_{\mathscr{M}} \gamma_{+}(T \boldsymbol{u}) \cdot \gamma_{+} \overline{\boldsymbol{v}} \mathrm{d} \sigma-\int_{\mathscr{M}} \gamma_{-}(T \boldsymbol{u}) \cdot \gamma_{-} \overline{\boldsymbol{v}} \mathrm{d} \sigma \tag{B.2.1}
\end{equation*}
$$

Under the symmetry hypothesis $\left(\mathfrak{H}_{\mathrm{A} 2}\right)$ we have the simplified second Green formula for all $\boldsymbol{u}, \boldsymbol{v} \in \mathrm{H}^{2}(\Omega)^{N}$

$$
\begin{align*}
\int_{\Omega}(\boldsymbol{u} \cdot \overline{L \boldsymbol{v}}-L \boldsymbol{u} \cdot \overline{\boldsymbol{v}}) \mathrm{d} \boldsymbol{y} & =\int_{\mathscr{M}}\left(\gamma_{+} \boldsymbol{u} \cdot \gamma_{+}(\overline{T \boldsymbol{v}})-\gamma_{-} \boldsymbol{u} \cdot \gamma_{-}(\overline{T \boldsymbol{v}})\right. \\
& \left.-\gamma_{+}(T \boldsymbol{u}) \cdot \gamma_{+} \overline{\boldsymbol{v}}+\gamma_{-}(T \boldsymbol{u}) \cdot \gamma_{-} \overline{\boldsymbol{v}}\right) \mathrm{d} \sigma \tag{B.2.2}
\end{align*}
$$

Let us recall a construction for the fundamental matrix of the operator $L\left(D_{\mathscr{X}}\right)$, i.e. the distribution $F_{L}$ such that

$$
\begin{equation*}
\forall \boldsymbol{x} \in \mathbb{R}^{n+1}, \quad L\left(D_{\mathscr{x}}\right) F_{L}(\boldsymbol{x})=\delta(\boldsymbol{x}) \mathrm{Id}, \quad F_{L} \in \mathscr{S}^{\prime}\left(\mathbb{R}^{n+1}\right) \tag{B.2.3}
\end{equation*}
$$

where Id is the identity matrix and $\delta$ is the Dirac distribution at 0

$$
\forall \varphi \in \mathscr{C}^{\infty}\left(\mathbb{R}^{n+1}\right), \quad(\delta, \varphi)=\varphi(0)
$$

After choosing in $\mathbb{R}^{n+1}$ a system of coordinates $\boldsymbol{x}=\left(x, x_{n+1}\right) \in \mathbb{R}^{n} \times \mathbb{R}$ which particularizes one coordinate, the fundamental matrix of equation (B.2.3) can be written in the following form, see [ Hr 1$]$ :

$$
\begin{equation*}
F_{L}(\boldsymbol{x}):=\mathscr{F}_{\xi \rightarrow x}^{-1}\left[\frac{1}{2 \pi} \int_{\mathscr{L}_{ \pm}} L^{-1}(\xi, \tau) e^{-i \tau x_{n+1}} \mathrm{~d} \tau\right] \quad \text { if } \quad \mp x_{n+1}>0 \tag{B.2.4}
\end{equation*}
$$

where $(\xi, \tau) \in \mathbb{R}^{n} \times \mathbb{R}$ represents the dual variables of $\left(x, x_{n+1}\right)$. The contour $\mathscr{L}_{+}\left(\mathscr{L}_{-}\right)$ is situated in the upper (in the lower) complex half-plane $\mathbb{C}_{+}:=\mathbb{R} \oplus i \mathbb{R}_{+}$(in $\mathbb{C}_{-}:=$ $\mathbb{R} \oplus i \mathbb{R}_{-}$) and is oriented counterclockwise (clockwise, respectively) encircling all roots of the polynomial $\operatorname{det} L(\xi, \tau)$ with respect to the variable $\tau$ in the corresponding half-planes $\tau \in \mathbb{C}_{ \pm}$.

Taking as test functions $\boldsymbol{v}(\boldsymbol{x})$ the columns of the matrix $F_{L}(\boldsymbol{x}-\boldsymbol{y})$ and inserting the equation $L\left(D_{\boldsymbol{x}}\right) \boldsymbol{u}=\boldsymbol{f}$ into the second Green formula (B.2.2), we easily obtain a representation formula for any $\boldsymbol{u}$ satisfying the equation $L\left(D_{\boldsymbol{x}}\right) \boldsymbol{u}=\boldsymbol{f}$ :

$$
\begin{equation*}
\forall \boldsymbol{x} \in \Omega, \quad \boldsymbol{u}(\boldsymbol{x})=N \boldsymbol{f}(\boldsymbol{x})+\mathscr{D}[\boldsymbol{u}](\boldsymbol{x})-\mathscr{V}[T \boldsymbol{u}](\boldsymbol{x}) \tag{B.2.5}
\end{equation*}
$$

where

$$
\forall \mathscr{X} \in \mathscr{M}, \quad[\boldsymbol{u}](\mathscr{X}):=\gamma_{+} \boldsymbol{u}(\mathscr{X})-\gamma_{-} \boldsymbol{u}(\mathscr{X}), \quad[T \boldsymbol{u}](\mathscr{X}):=\gamma_{+} T \boldsymbol{u}(\mathscr{X})-\gamma_{-} T \boldsymbol{u}(\mathscr{X})
$$

denote the jumps of the functions $\boldsymbol{u}(\mathscr{X})$ and $T \boldsymbol{u}(x)$ across the surface $\mathscr{M}$; the operators $\mathscr{V}, \mathscr{D}$ and $N$ are the well-known single layer, double layer and volume (Newton) potentials:

$$
\begin{gather*}
\mathscr{V} \boldsymbol{\phi}(\boldsymbol{x})=\int_{\mathscr{M}} F_{L}(\boldsymbol{x}-\sigma) \boldsymbol{\phi}(\sigma) \mathrm{d} \sigma, \quad \mathscr{D} \boldsymbol{\phi}(\boldsymbol{x})=\int_{\mathscr{M}}\left(T F_{L}\right)^{*}(\sigma-\boldsymbol{x}) \boldsymbol{\phi}(\sigma) \mathrm{d} \sigma,  \tag{B.2.6}\\
N \boldsymbol{f}(\boldsymbol{x})=\int_{\mathbb{R}^{n+1}}\left(F_{L}\right)^{*}(\boldsymbol{x}-\boldsymbol{y}) \boldsymbol{f}(\boldsymbol{y}) \mathrm{d} \boldsymbol{y}, \quad \boldsymbol{x} \in \Omega . \tag{B.2.7}
\end{gather*}
$$

Here $\mathscr{A}^{*}:=\overline{\mathscr{A}^{\top}}$ denotes the hermitian conjugate of the matrix $\mathscr{A}$.
Solving the boundary value problem (A.1.1) with the help of the representation formula (B.2.5) we have to find only one density, either $\boldsymbol{\varphi}=[\boldsymbol{u}] \in \widetilde{\mathrm{H}}^{\frac{1}{2}}(\mathscr{M})$ for the Neumann problem or $\boldsymbol{\psi}=[T \boldsymbol{u}] \in \widetilde{\mathrm{H}}^{-\frac{1}{2}}(\mathscr{M})$ for the Dirichlet problem (due to the boundary conditions in (B.2.5) the other density vanishes on $\mathscr{M}$ ). Invoking the well-known jump relations ("Plemelj formulae") (see, e.g., $[\mathrm{KuGeBaBu} 1, \mathrm{ChPi}]$ ) we get the following pseudodifferential equations on the crack surface (compare with (A.3.1))

$$
\begin{array}{rlrl}
W\left(x, D_{x}\right) \boldsymbol{\varphi}(x) & =-\gamma_{+} T N \boldsymbol{f}(x), & x \in \mathscr{M}, & \\
\text { for Neumann, }  \tag{B.2.9}\\
V\left(x, D_{\mathscr{x}}\right) \boldsymbol{\psi}(x)=\gamma_{+} N \boldsymbol{f}(x), & \mathscr{x} \in \mathscr{M}, & & \text { for Dirichlet. }
\end{array}
$$

Here $W\left(\mathscr{x}, D_{\mathscr{E}}\right)=\gamma_{+} T \mathscr{D}=\gamma_{-} T \mathscr{D}$ is the trace of the composition of the Neumann operator with the double layer potential and is a hypersingular operator, understood as a pseudodifferential operator of order 1. $V\left(\mathscr{x}, D_{\mathscr{X}}\right)=\gamma_{+} \mathscr{V}=\gamma_{-} \mathscr{V}$ is the trace of the single layer potential on the surface $\mathscr{M}$ and is a weakly singular integral operator (pseudodifferential operator of order -1 ).

Thus, by solving the equation (B.2.8) or (B.2.9), and inserting the solution into the representation formula

$$
\begin{array}{lll}
\boldsymbol{u}(\boldsymbol{x})=N \boldsymbol{f}(\boldsymbol{x})+\mathscr{D} \boldsymbol{\varphi}(\boldsymbol{x}), & \boldsymbol{x} \in \Omega, & \text { for Neumann } \\
\boldsymbol{u}(\boldsymbol{x})=N \boldsymbol{f}(\boldsymbol{x})-\mathscr{V} \boldsymbol{\psi}(\boldsymbol{x}), & \boldsymbol{x} \in \Omega, & \text { for Dirichlet, } \tag{B.2.11}
\end{array}
$$

we obtain a solution of the boundary value problem (A.1.1).

## B. 3 ASYMPTOTICS OF SOLUTIONS OF $\Psi$ DE - A GENERAL CASE

In this section we recall general results on asymptotics of solutions to $\Psi D E$ on a manifold with smooth boundary from [Es1, ChkDu1] obtained by the Wiener-Hopf approach.

Let us consider a classical $N \times N$ matrix symbol $\mathbf{a}(\mathscr{X} ; \xi)$ of order $\nu \in \mathbb{R}$, defined on the cotangent manifold $\mathscr{T}^{*} \mathscr{M}$ to $\mathscr{M} \subset \mathbb{R}^{n+1}$ :

$$
\begin{gather*}
\mathbf{a} \in S_{\mathrm{cl}}^{\nu}\left(\mathscr{T}^{*} \mathscr{M}\right)^{N \times N} \Longleftrightarrow \mathbf{a}(\mathscr{X}, \xi)=\mathbf{a}_{0}(\mathscr{X}, \xi)+\mathbf{a}_{1}(\mathscr{X}, \xi)+\cdots, \\
\forall \lambda>0, \forall \mathscr{X} \in \mathscr{M}, \forall \xi \in \mathbb{R}^{n}, \quad \mathbf{a}_{j}(x, \lambda \xi)=\lambda^{\nu-j} \mathbf{a}_{j}(x, \xi), \tag{B.3.1}
\end{gather*}
$$

where $\mathbf{a}_{j}(x, \theta)$ are $\mathscr{C}^{\infty}$-smooth on the bundle of cotangent unit spheres $\overline{\mathscr{M}} \times \mathbb{S}^{n-1} \subset$ $\mathscr{T}^{*} \mathscr{M}$ (see [Hr1] and [ChkDu1, § 1.2]).

For any Sobolev exponent $s \in \mathbb{R}$, the corresponding $N \times N$ system of $\Psi \mathrm{DE}$ on $\mathscr{M}$ with symbol $\mathbf{a}(\mathscr{x} ; \xi)$ is continuous from $\widetilde{\mathrm{H}}^{s}(\mathscr{M})^{N}$ into $\mathrm{H}^{s-\nu}(\mathscr{M})^{N}$. We are interested in the structure of any $\phi$ satisfying for some $s \in \mathbb{R}$ and an integer $K>0$ :

$$
\begin{equation*}
\phi \in \widetilde{\mathrm{H}}^{s}(\mathscr{M})^{N} \quad \text { such that } \quad \mathbf{a}\left(\mathscr{x} ; D_{\mathscr{X}}\right) \boldsymbol{\phi}=\boldsymbol{g}, \quad \text { with } \quad \boldsymbol{g} \in \mathrm{H}^{s-\nu+K}(\mathscr{M})^{N} . \tag{B.3.2}
\end{equation*}
$$

Further we suppose that the principal homogeneous part $\mathbf{a}_{0}(x ; \xi)$, which we will also denote by $\mathbf{a}_{\mathrm{pr}}(x ; \xi)$ is elliptic, which reads

$$
\begin{equation*}
\operatorname{det} \mathbf{a}_{\mathrm{pr}}(\mathscr{X} ; \xi) \neq 0, \quad \mathscr{x} \in \overline{\mathscr{M}}, \quad \xi \in \mathbb{R}^{n} \backslash\{0\} \tag{B1}
\end{equation*}
$$

The following $N \times N$ matrix plays a fundamental role in the structure of the solutions $\phi$ satisfying (B.3.2)

$$
\begin{equation*}
\mathbf{b}\left(x^{\prime}\right):=\left[\mathbf{a p r}_{\mathrm{pr}}\left(x^{\prime}, 0 ; 0,+1\right)\right]^{-1} \mathbf{a}_{\mathrm{pr}}\left(x^{\prime}, 0 ; 0,-1\right), \quad x^{\prime} \in \mathscr{E} \tag{B.3.3}
\end{equation*}
$$

where we recall that $\mathscr{x}:=\left(\mathscr{X}^{\prime}, \mathscr{X}_{n}\right) \in \mathscr{M}$ are the local and $\xi=\left(\xi^{\prime}, \xi_{n}\right)$ are the dual coordinates, with $\mathscr{X}^{\prime} \in \mathscr{E}=\partial \mathscr{M}$ the edge variable. Note that for all $\xi^{\prime} \in \mathbb{R}^{n-1}$ :

$$
\mathbf{a}_{\mathrm{pr}}\left(x^{\prime}, 0 ; 0, \pm 1\right)=\lim _{t \rightarrow \pm \infty}|t|^{-\nu} \mathbf{a}_{\mathrm{pr}}\left(x^{\prime}, 0 ; \xi^{\prime}, t\right)
$$

For any $\mathscr{x}^{\prime} \in \mathscr{E}$, let us denote by

$$
\lambda_{1}\left(\mathscr{X}^{\prime}\right), \ldots, \lambda_{N}\left(\mathscr{X}^{\prime}\right) \text { the eigenvalues of } \mathbf{b}\left(\mathscr{X}^{\prime}\right),
$$

where we repeat each eigenvalue according to its algebraic multiplicity.
The assumption which will ensure the absence of logarithms in the principal term of the asymptotics of the $\phi$ satisfying (B.3.2) is that $\mathbf{b}$ is diagonalizable in each point $\mathscr{X}^{\prime}$ in $\mathscr{E}$, and that the eigenvalues are $\mathscr{C}{ }^{\infty}(\mathscr{E})$, which is written as:
$\left(\mathfrak{H}_{\mathrm{B} 2}\right) \quad\left\{\begin{array}{l}\forall \mathscr{X}^{\prime} \in \mathscr{E}, \exists \text { a numbering of eigenvalues and an invertible matrix } \mathscr{K}\left(x^{\prime}\right): \\ \mathbf{b}\left(\mathscr{X}^{\prime}\right)=\mathscr{K}\left(\mathscr{X}^{\prime}\right)\left(\operatorname{diag}\left\{\lambda_{1}\left(\mathscr{X}^{\prime}\right), \ldots, \lambda_{N}\left(\mathscr{X}^{\prime}\right)\right\}\right) \mathscr{K}^{-1}\left(x^{\prime}\right) \\ \mathscr{X}^{\prime} \mapsto \mathscr{K}\left(\mathscr{X}^{\prime}\right), \quad \mathscr{X}^{\prime} \mapsto \lambda_{j}\left(\mathscr{X}^{\prime}\right) \text { are } \mathscr{C}^{\infty}(\mathscr{E}) .\end{array}\right.$
We need one more assumption on the eigenvalues of $\mathbf{b}\left(x^{\prime}\right)$ : let us set

$$
\delta_{j}\left(x^{\prime}\right)=(2 \pi i)^{-1} \log \lambda_{j}\left(x^{\prime}\right), \quad j=1, \ldots, N .
$$

We assume that
$\left(\mathfrak{H}_{\mathrm{B} 3}\right) \quad\left\{\begin{array}{l}\exists \eta \in\left(-\frac{1}{2}, \frac{1}{2}\right), \quad \exists \mathrm{a} \mathscr{C}^{\infty}(\mathscr{E}) \text { determination of the } \delta_{j}\left(\mathscr{X}^{\prime}\right) \\ \text { such that } \quad \forall \mathscr{X}^{\prime} \in \mathscr{E}, \quad \eta-\frac{1}{2}<\operatorname{Re} \delta_{j}\left(\mathscr{X}^{\prime}\right)<\eta+\frac{1}{2} .\end{array}\right.$
While locally a consequence of $\left(\mathfrak{H}_{\mathrm{B} 2}\right)$, this assumption has to be required to hold globally on $\mathscr{E}$.

The following result, see [DuSaWe1, Lemma A.6], provides a general framework where assumptions $\left(\mathfrak{H}_{\mathrm{B} 2}\right)$ and $\left(\mathfrak{H}_{\mathrm{B} 3}\right)$ are satisfied.

Lemma B.3.1 Iffor any $\mathscr{X}^{\prime} \in \mathscr{E}$ the two matrices $\mathrm{a}_{\mathrm{pr}}\left(\mathscr{X}^{\prime}, 0 ; 0, \pm 1\right)$ in (B.3.3) are positive definite, then the matrix $\mathbf{b}\left(\mathscr{X}^{\prime}\right)$ is diagonalizable with unitary $\mathscr{K}\left(\mathscr{X}^{\prime}\right)$, its eigenvalues are real, which means that the numbers $\delta_{j}\left(\mathscr{X}^{\prime}\right)$ can be chosen purely imaginary:

$$
\begin{equation*}
\delta_{j} \in \mathscr{C}^{\infty}(\mathscr{E}), \quad \operatorname{Re} \delta_{j}\left(\mathscr{X}^{\prime}\right)=0 \quad \text { for all } \quad j=1, \cdots, N . \tag{B.3.4}
\end{equation*}
$$

The main result in this section is the asymptotic structure of solutions $\phi$ of (B.3.2), whose first term does not contain logarithms. We recall that $\left(\mathscr{X}^{\prime}, r\right)=\left(\mathscr{X}^{\prime}, r, \pm \pi\right)$ denotes the local cylindrical coordinate system on $\mathscr{M}$ in a closed tubular neighborhood of the edge $\mathscr{E}=\partial \mathscr{M}$ (see Definition A.4.1).

Theorem B.3.2 (see [ChkDu1] and [Es1, Ch.26]). We assume hypotheses ( $\mathfrak{H}_{\mathrm{B} 1}$ ), ( $\mathfrak{H}_{\mathrm{B} 2}$ ) and $\left(\mathfrak{H}_{\mathrm{B} 3}\right)$. We choose

- a determination of the $\delta_{j}, j=1, \ldots, N$,
- a real Sobolev exponent $s$,
such that there holds for all $\mathscr{X}^{\prime} \in \mathscr{E}$

$$
\begin{equation*}
\forall j=1, \ldots, N, \quad \frac{\nu}{2}+\operatorname{Re} \delta_{j}\left(\mathscr{X}^{\prime}\right)>-1 \tag{B.3.5}
\end{equation*}
$$

and

$$
\begin{equation*}
\forall j=1, \ldots, N, \quad-\frac{\nu}{2}+s-\frac{1}{2}<\operatorname{Re} \delta_{j}\left(\mathscr{X}^{\prime}\right)<-\frac{\nu}{2}+s+\frac{1}{2} . \tag{B.3.6}
\end{equation*}
$$

Moreover let $\phi \in \widetilde{\mathrm{H}}^{s}(\mathscr{M})^{N}$ be a solution of the equation $\mathbf{a}\left(\mathscr{x} ; D_{\mathscr{X}}\right) \boldsymbol{\phi}=\boldsymbol{g}$ where the right hand side $\boldsymbol{g}$ is $\mathscr{C}^{\infty}(\mathscr{M})^{N}$. Then, for any integer $K>0$ the solution $\phi$ has the following asymptotic expansion

$$
\begin{array}{r}
\boldsymbol{\phi}\left(\mathscr{X}^{\prime}, r\right)=\mathscr{K}\left(\mathscr{X}^{\prime}\right) r^{\frac{\nu}{2}+\Delta\left(\mathscr{X}^{\prime}\right)} \chi(r)\left[\boldsymbol{d}^{0}\left(\mathscr{X}^{\prime}\right)+\sum_{k=1}^{K-1} r^{k} \sum_{q=0}^{\sigma(k)} \boldsymbol{d}^{k, q}\left(\mathscr{X}^{\prime}\right) \log ^{q} r\right] \\
+\phi_{\mathrm{rem}, K}\left(\mathscr{X}^{\prime}, r\right), \quad \phi_{\mathrm{rem}, K} \in \widetilde{\mathrm{H}}^{s+K}(\mathscr{M})^{N} \tag{B.3.7}
\end{array}
$$

with $N$-vector coefficients $\boldsymbol{d}^{0}, \boldsymbol{d}^{k, q}$ in $\mathscr{C}^{\infty}(\mathscr{E})^{N}$. Here, the vector $\Delta$ is defined as $\left(\delta_{1}, \cdots, \delta_{N}\right)^{\top}$ and for any $\mu \in \mathbb{R}, r^{\mu+\Delta}$ is understood as the diagonal $N \times N$ matrix

$$
\begin{equation*}
r^{\mu+\Delta}:=\operatorname{diag}\left\{r^{\mu+\delta_{1}}, \ldots, r^{\mu+\delta_{N}}\right\} \tag{B.3.8}
\end{equation*}
$$

Remark B.3.3 (i) In [Es1, Ch.26], it is proved that the asymptotics of $\phi$ has no logarithmic term in its leading summand, and in [ChkDu1] the more explicit formula (B.3.7) is proved.
(ii) It is possible to extend hypothesis $\left(\mathfrak{H}_{\mathrm{B} 2}\right)$ to certain cases where $\mathbf{b}\left(\mathscr{X}^{\prime}\right)$ is not diagonalizable: then we assume that we have a canonical Jordan decomposition with a $\mathscr{C}^{\infty}(\mathscr{E})$ dependence. This implies in particular that the geometrical multiplicities are constant along $\mathscr{E}$. Then it is proved in [ChkDu1] that there holds a decomposition like (B.3.7), with explicit
logarithmic terms in the leading summand of the asymptotics. This means that the condition $\left(\mathfrak{H}_{\mathrm{B} 2}\right)$ is necessary and sufficient so that logarithms are absent in the leading summand of the asymptotic of a solution (B.3.7).
(iii) It is possible to get the first term of the asymptotic expansion without the smoothness properties on $\mathscr{K}$ and $\delta_{j}$, but the further terms are not available so far, see [ChkDu1].

## B. 4 ASYMPTOTICS OF $\Psi$ DE - SYMBOLS WITH CONTINUITY PROPERTY

Here are the conditions which ensure that logarithms disappear from the entire asymptotics (B.3.7). These conditions apply to the full symbol $\sum_{j \geq 0} \mathbf{a}_{j}\left(\mathscr{X}^{\prime}, \mathscr{X}_{n} ; \xi^{\prime}, \xi_{n}\right)$ :
$\left(\mathfrak{H}_{\mathrm{B} 4}\right) \quad\left\{\begin{array}{l}\forall \mathscr{X}^{\prime} \in \mathscr{E}, \quad \forall j \in \mathbb{N}_{0}, \quad \forall \alpha^{\prime} \in \mathbb{N}_{0}^{n-1}, \quad m \in \mathbb{N}_{0}, \\ \left(\partial_{\mathscr{X}_{n}}^{m} \partial_{\xi^{\prime}}^{\alpha^{\prime}} \mathbf{a}_{j}\right)\left(\mathscr{X}^{\prime}, 0 ; 0,-1\right)=(-1)^{j+\left|\alpha^{\prime}\right|}\left(\partial_{\mathscr{X}_{n}}^{m} \partial_{\xi^{\prime}}^{\alpha^{\prime}} \mathbf{a}_{j}\right)\left(\mathscr{X}^{\prime}, 0 ; 0,+1\right) .\end{array}\right.$
We note that the above condition implies that

$$
\begin{equation*}
\forall \beta \in \mathbb{N}_{0}^{n}, \quad\left(\partial_{\mathscr{X}}^{\beta} \partial_{\xi^{\prime}}^{\alpha^{\prime}} \mathbf{a}_{j}\right)\left(\mathscr{x}^{\prime}, 0 ; 0,-1\right)=(-1)^{j+\left|\alpha^{\prime}\right|}\left(\partial_{\mathscr{X}}^{\beta} \partial_{\xi^{\prime}}^{\alpha^{\prime}} \mathbf{a}_{j}\right)\left(\mathscr{X}^{\prime}, 0 ; 0,+1\right) \tag{B.4.1}
\end{equation*}
$$

On the other hand, concerning the principal symbol, the above condition implies that

$$
\mathbf{a}_{\mathrm{pr}}\left(\mathscr{X}^{\prime}, 0 ; 0,-1\right)=\mathbf{a}_{\mathrm{pr}}\left(\mathscr{X}^{\prime}, 0 ; 0,+1\right)
$$

whence for all $x^{\prime} \in \mathscr{E}, \mathbf{b}\left(\mathscr{X}^{\prime}\right)=\mathrm{Id}$. Thus condition $\left(\mathfrak{H}_{\mathrm{B} 4}\right)$ implies conditions $\left(\mathfrak{H}_{\mathrm{B} 2}\right)$ and $\left(\mathfrak{H}_{\mathrm{B} 3}\right)$.

The main result about asymptotics without logarithmic terms within the Wiener-Hopf approach are formulated in Theorem B.4.1.
Theorem B.4.1 Let $\mathbf{a}(\mathscr{x} ; \xi)$ be a classical symbol (B.3.1) of order $\nu>-2$ and let its homogeneous components $\mathbf{a}_{0}(\mathscr{x} ; \xi), \mathbf{a}_{1}(\mathscr{X} ; \xi), \ldots$ satisfy the continuity property $\left(\mathfrak{H}_{\mathrm{B} 4}\right)$ on the boundary $\mathscr{E}$. Let s be a Sobolev exponent such that

$$
\begin{equation*}
\frac{\nu}{2}-\frac{1}{2}<s<\frac{\nu}{2}+\frac{1}{2} \tag{B.4.2}
\end{equation*}
$$

Let $\phi \in \widetilde{\mathrm{H}}^{s}(\mathscr{M})^{N}$ be a solution of the equation $\mathbf{a}\left(\mathscr{x} ; D_{\mathscr{x}}\right) \boldsymbol{\phi}=\boldsymbol{g}$ where the right hand side $\boldsymbol{g}$ is $\mathscr{C}^{\infty}(\mathscr{M})^{N(1)}$, the solution has the following asymptotic expansion for any integer $K>0$

$$
\begin{equation*}
\phi=\sum_{k=0}^{K-1} r^{\frac{\nu}{2}+k} \chi(r) \boldsymbol{d}^{k}\left(\mathscr{X}^{\prime}\right)+\boldsymbol{\phi}_{\mathrm{rem}, K}, \quad \phi_{\mathrm{rem}, K} \in \widetilde{\mathrm{H}}^{s+K}(\mathscr{M})^{N}, \tag{B.4.3}
\end{equation*}
$$

where the $N$-vectors $\boldsymbol{d}^{k}, k=0,1, \ldots$ belong to $\mathscr{C}^{\infty}(\mathscr{E})$.

[^0]We postpone the proof of the theorem until § B.7.
The assumptions of Theorem B.4.1 hold for the boundary $\Psi$ DE (B.2.8) and (B.2.9) corresponding to the BVP (A.1.1). This follows from the following theorem.

Theorem B.4.2 The symbols of the boundary $\Psi D E$ (B.2.8) of order $\nu=1$ and (B.2.9) of order $\nu=-1$ are positive definite and satisfy the continuity property $\left(\mathfrak{H}_{\mathrm{B} 4}\right)$. Moreover, for any volume data $f \in \mathscr{C}_{0}^{\infty}\left(\mathbb{R}^{n+1}\right)$, the right hand sides of equations (B.2.8) and (B.2.9) are in $\mathscr{C}^{\infty}(\mathscr{M})^{N}$, and equations (B.2.8) and (B.2.9) have unique solutions $\varphi \in \widetilde{\mathrm{H}}^{s}(\mathscr{M})^{N}$ and $\psi \in \widetilde{\mathrm{H}}^{s-1}(\mathscr{M})^{N}$, respectively, for any $s \in(0,1)$. Thus asymptotics (B.4.3) hold for these solutions.

Proof. We quote [ChkDu1, CoSt1, DuNaSh1, DuWe1] for the proofs of positive definiteness of the symbols and unique solvability (also in more general spaces) of $\Psi D E$ (B.2.8) and (B.2.9) and concentrate on the proof of the continuity property (B.4.1).

In [ChkDu1, Example 1.17] it is proved that the symbols of both equations (B.2.8) and (B.2.9) are classical ${ }^{(2)}$ and the components of the asymptotic representation of the symbols have the following form

$$
\begin{aligned}
W(x ; \xi) & =W_{0}(x ; \xi)+W_{1}(x ; \xi)+\cdots+W_{j}(x ; \xi)+\cdots \\
V(x ; \xi) & =V_{0}(x ; \xi)+V_{1}(x ; \xi)+\cdots+V_{j}(x ; \xi)+\cdots
\end{aligned}
$$

where the homogeneous components $W_{j}(x ; \xi)$ and $V_{j}(x ; \xi)$ (of orders $1-j$ and $-1-j$ respectively) are generated by an explicit symbol $w$ and $v$ respectively

$$
\begin{align*}
W_{j}(x ; \xi) & =\sum_{\substack{|\alpha|-|\beta| \leq j \geq 0 \\
2 \beta \leq \alpha}} \mathbf{a}_{\alpha, \beta}(x) \xi^{\beta} \partial_{\xi}^{\alpha} \boldsymbol{w}(x ; \xi),  \tag{B.4.4}\\
V_{j}(x ; \xi) & =\sum_{\substack{|\alpha|-|\beta|=j \geq 0 \\
2 \beta \leq \alpha}} \mathbf{a}_{\alpha, \beta}(x) \xi^{\beta} \partial_{\xi}^{\alpha} \boldsymbol{v}(x ; \xi), \tag{B.4.5}
\end{align*}
$$

where the sums are finite since $|\alpha|-|\beta|=j$ and $2 \beta \leq \alpha$ imply $2|\beta| \leq|\alpha| \leq 2 n j$, and where the matrices $\mathbf{a}_{\alpha, \beta}(x)$ have $\mathscr{C}^{\infty}(\mathscr{M})$ coefficients.

The generating symbols $w, v$ are defined for $\mathscr{x} \in \mathscr{M}$ and $\xi \in \mathbb{R}^{n}$ as follows - the contour $\mathscr{L}_{+}$is the same as in (B.2.4) and the Jacobian $\mathscr{J}_{\kappa}\left(\mathscr{X}^{\prime}\right)$ as in Definition A.4.1 (vi):

$$
\begin{align*}
& \boldsymbol{w}(\mathscr{X} ; \xi):=\int_{\mathscr{L}_{+}} T\left(\mathscr{X} ; \mathscr{F}_{\kappa}(\mathscr{x})(\xi, \tau)\right) L^{-1}\left(\mathscr{J}_{\kappa}(\mathscr{x})(\xi, \tau)\right)^{\top}  \tag{B.4.6}\\
& T\left(\mathscr{x} ; \mathscr{J}_{\kappa}(\mathscr{x})(\xi, \tau)\right)^{\top} \mathrm{d} \tau \\
& \boldsymbol{v}(\mathscr{X} ; \xi):=\int_{-\infty}^{\infty} L^{-1}\left(\mathscr{J}_{\kappa}(\mathscr{x})(\xi, \tau)\right) \mathrm{d} \tau . \tag{B.4.7}
\end{align*}
$$

[^1]In particular, the principal symbols $W_{\mathrm{pr}}(x ; \xi)=W_{0}(x ; \xi)$ and $V_{\mathrm{pr}}(x ; \xi)=V_{0}(x ; \xi)$ both have the following coefficient

$$
\mathbf{a}_{0,0}(x):=\frac{\Gamma_{\varkappa}(x)}{2 \pi \operatorname{det} D \varkappa(x)},
$$

where $\Gamma_{\varkappa}(\mathscr{x})$ is the Gram determinant of the local coordinate diffeomorphisms $\varkappa$.
Since the elliptic differential operator $L\left(D_{\boldsymbol{x}}\right)$ in (A.1.1) is supposed to be homogeneous of degree 2 , its symbol $L\left(\xi, \xi_{n+1}\right)$ is even

$$
\forall \boldsymbol{\xi}=\left(\xi, \xi_{n+1}\right) \in \mathbb{R}^{n+1}, \quad L\left(-\xi,-\xi_{n+1}\right)=L\left(\xi, \xi_{n+1}\right) .
$$

As a consequence, with the change of variable $\tau \mapsto-\tau$ in integrals (B.4.6) and (B.4.7), we find that the generating symbols $V$ and $w$ are even ${ }^{(3)}$

$$
\forall x \in \mathscr{M}, \quad \forall \xi \in \mathbb{R}^{n}, \quad v(x,-\xi)=v(\mathscr{x}, \xi) \quad \text { and } \quad \boldsymbol{w}(\mathscr{x},-\xi)=\boldsymbol{w}(\mathscr{x}, \xi) .
$$

Therefore, as a consequence of formulas (B.4.6) and (B.4.7), for all $\mathscr{X} \in \mathscr{M}$, for all $\lambda \in \mathbb{R}$, for all integers $j, m=0,1, \ldots$ and all multiindices $\alpha^{\prime} \in \mathbb{N}_{0}^{n-1}$, there holds

$$
\begin{aligned}
\left(\partial_{\mathscr{x}_{n}}^{m} \partial_{\xi^{\prime}}^{\alpha^{\prime}} W_{j}\right)\left(\mathscr{X}^{\prime}, 0 ; 0,-\lambda\right) & =(-1)^{j+\left|\alpha^{\prime}\right|}\left(\partial_{x_{n}}^{m} \partial_{\xi^{\prime}}^{\alpha^{\prime}} W_{j}\right)\left(\mathscr{X}^{\prime}, 0 ; 0, \lambda\right), \\
\left(\partial_{x_{n}}^{m} \partial_{\xi^{\prime}}^{\alpha^{\prime}} V_{j}\right)\left(\mathscr{x}^{\prime}, 0 ; 0,-\lambda\right) & =(-1)^{j+\left|\alpha^{\prime}\right|}\left(\partial_{x_{n}}^{m} \partial_{\xi^{\prime}}^{\alpha^{\prime}} V_{j}\right)\left(\mathscr{X}^{\prime}, 0 ; 0, \lambda\right) .
\end{aligned}
$$

## B. 5 TDE IN DIMENSION 1

Before we start the proof of the main theorem B.4.1, we want to explain the principal mechanism responsible for the absence of logarithmic terms by presenting the result in a very simple situation, namely the case of a scalar elliptic pseudodifferential equation with constant coefficients on the half-line $\mathbb{R}_{+}$. This simple one-dimensional situation allows us to stay free of many of the technical difficulties of the higher-dimensional case and to concentrate on the essential feature, namely the role of the continuity condition for the asymptotic expansion of the symbol. We can show in this case that a natural generalization of this condition is not only sufficient, but also necessary for the absence of logarithmic terms in the asymptotics of the solution. The class of operators considered here can be larger than the one obtained from the 2 D crack problem.

We need the following well-known Fourier transform of distributions supported in the positive half-line, see for instance [Es1]. By $\chi_{+}$and $\chi_{-}$we denote the characteristic functions of $\mathbb{R}_{+}$and $\mathbb{R}_{-}$, respectively.

Lemma B.5.1
(i) $\mathscr{F}_{t \rightarrow \lambda}\left(\chi_{+}(t) t^{\mu-1} e^{-\tau t}\right)=\Gamma(\mu) e^{i \frac{\pi}{2} \mu}(\lambda+i \tau)^{-\mu}, \quad \tau>0$

[^2](ii) $\mathscr{F}_{t \rightarrow \lambda}\left(\chi_{+}(t) \log t t^{\mu-1} e^{-\tau t}\right)=(\lambda+i \tau)^{-\mu}(c \log (\lambda+i \tau)+d)$ with $c=-\Gamma(\mu) e^{i \frac{\pi}{2} \mu}$ and $d=\frac{d}{d \mu}\left(\Gamma(\mu) e^{i \frac{\pi}{2} \mu}\right)$.
Another crucial result concerns the additive decomposition of homogeneous distributions into "plus" and "minus" terms.

Lemma B.5. 2 Let $a^{+}, a^{-}, \gamma \in \mathbb{C}$. Then
(i) If $\gamma \notin \mathbb{Z}$, we have the representation

$$
\left(a^{+} \chi_{+}(t)+a^{-} \chi_{-}(t)\right)|t|^{\gamma}=\frac{a^{-}-e^{-i \pi \gamma} a^{+}}{e^{i \pi \gamma}-e^{-i \pi \gamma}}(t+i 0)^{\gamma}-\frac{a^{-}-e^{i \pi \gamma} a^{+}}{e^{i \pi \gamma}-e^{-i \pi \gamma}}(t-i 0)^{\gamma}
$$

(ii) $\gamma \in \mathbb{Z}$, we have the representation

$$
\begin{aligned}
\left(a^{+} \chi_{+}(t)+a^{-} \chi_{-}(t)\right)|t|^{\gamma}=a^{+}(t+i 0)^{\gamma}+\frac{(-1)^{\gamma} a^{-}-a^{+}}{2 i \pi} & \left((t+i 0)^{\gamma} \log (t+i 0)\right. \\
& \left.-(t-i 0)^{\gamma} \log (t-i 0)\right)
\end{aligned}
$$

Proof. It suffices to use the identities

$$
\begin{aligned}
(t \pm i 0)^{\gamma} & =\chi_{+}(t) t^{\gamma}+\chi_{-}(t) e^{ \pm i \pi \gamma}|t|^{\gamma} \\
\log (t \pm i 0) & =\chi_{+}(t) \log t+\chi_{-}(t)(\log |t| \pm i \pi)
\end{aligned}
$$

Let $a \in \mathscr{C}^{\infty}(\mathbb{R} ; \mathbb{C})$ be a classical elliptic symbol of order $\nu \in \mathbb{R}$ with constant coefficients, i.e. $a(\xi) \neq 0$ for all $\xi \in \mathbb{R}$, and $a$ has an asymptotic expansion in homogeneous terms

$$
\begin{equation*}
a(\xi) \sim \sum_{j=0}^{\infty} a_{j}(\xi) \quad \text { with } \quad \forall j \in \mathbb{N}, \forall t>0, \forall \xi \in \mathbb{R}, a_{j}(t \xi)=t^{\nu-j} a_{j}(\xi) \tag{B.5.1}
\end{equation*}
$$

In one dimension, homogeneous functions are determined by two values:

$$
\begin{equation*}
a_{j}(\xi)=\left(a_{j}^{+} \chi_{+}(\xi)+a_{j}^{-} \chi_{-}(\xi)\right)|\xi|^{\nu-j} \tag{B.5.2}
\end{equation*}
$$

From the ellipticity follows that $a_{0}^{+} a_{0}^{-} \neq 0$, and we can define

$$
\lambda=\frac{a_{0}^{-}}{a_{0}^{+}} \in \mathbb{C} .
$$

By $p_{+} a(D) u$ we denote the restriction of

$$
a(D) u(x):=\mathscr{F}_{\xi \mapsto x}^{-1} a(\xi)(\mathscr{F} u)(\xi)
$$

to the half-line $\mathbb{R}_{+}$.

Theorem B.5.3 Let a be a classical elliptic symbol of order $\nu$ as above, and let $\delta \in \mathbb{C}$, $s \in \mathbb{R}$ be chosen such that

$$
\begin{equation*}
e^{2 i \pi \delta}=\lambda \quad \text { and } \quad \frac{\nu}{2}+\operatorname{Re} \delta-\frac{1}{2}<s<\frac{\nu}{2}+\operatorname{Re} \delta+\frac{1}{2} \tag{B.5.3}
\end{equation*}
$$

and $\frac{\nu}{2}+\delta \notin\{-1,-2, \ldots\}$. Let $u \in \widetilde{\mathrm{H}}^{s}\left(\mathbb{R}_{+}\right)$be solution of

$$
\begin{equation*}
p_{+} a(D) u=g \quad \text { on } \quad \mathbb{R}_{+} \tag{B.5.4}
\end{equation*}
$$

with $g \in \mathscr{C}^{\infty}\left(\overline{\mathbb{R}}_{+}\right) \cap \mathrm{H}^{s-\nu}\left(\mathbb{R}_{+}\right)$. Then $u$ has an asymptotic expansion as $x \rightarrow 0$ :

$$
u(x) \sim \sum_{k \geq 0} \sum_{q=0}^{q_{k}} c_{k q} x^{\frac{\nu}{2}+\delta+k} \log ^{q} x .
$$

This asymptotic expansion for any such $u$ is free of logarithms, i.e. $q_{k}=0$ for all $k \geq 0$, if and only if the following condition $\left(\mathfrak{H}_{\mathrm{B} 5}\right)$ is satisfied

$$
\left(\mathfrak{H}_{\mathrm{B} 5}\right) \quad \forall j \geq 0, \quad a_{j}^{-}=(-1)^{j} \lambda a_{j}^{+} .
$$

Note that, reduced to the case of dimension 1 with scalar operators, condition $\left(\mathfrak{H}_{\mathrm{B} 5}\right)$ is a generalization of condition $\left(\mathfrak{H}_{\mathrm{B} 4}\right)$ which corresponds to taking $\lambda=1$.
Proof. (i) We first show the sufficiency of condition $\left(\mathfrak{H}_{\mathrm{B} 5}\right)$.
If $\left(\mathfrak{H}_{\mathrm{B} 5}\right)$ is satisfied, then we can write

$$
\begin{equation*}
a(\xi)=a_{0}(\xi) q(\xi), \tag{B.5.5}
\end{equation*}
$$

where $q(\xi)$ has an asymptotic expansion of the form

$$
\begin{equation*}
q(\xi) \sim 1+\sum_{j \geq 1} q_{j} \xi^{-j} \quad \text { with } \quad q_{j}=\frac{a_{j}^{+}}{a_{0}^{+}} \tag{B.5.6}
\end{equation*}
$$

Thus $q$ is a symbol of rational type.
For $a_{0}$ we find the factorization

$$
\begin{equation*}
a_{0}(\xi)=a_{0}^{+}(\xi+i 0)^{\frac{\nu}{2}+\delta}(\xi-i 0)^{\frac{\nu}{2}-\delta} \tag{B.5.7}
\end{equation*}
$$

Let us introduce the corresponding $\mathscr{C}^{\infty}(\mathbb{R})$ symbol

$$
\begin{equation*}
a^{\infty}(\xi)=a_{0}^{+}(\xi+i)^{\frac{\nu}{2}+\delta}(\xi-i)^{\frac{\nu}{2}-\delta} \tag{B.5.8}
\end{equation*}
$$

Then we have the global representation of the symbol $a$ as the product

$$
\begin{equation*}
a(\xi)=a^{\infty}(\xi) q^{\infty}(\xi) \tag{B.5.9}
\end{equation*}
$$

with a symbol of rational type

$$
\begin{equation*}
q^{\infty}(\xi) \sim 1+\sum_{j \geq 1} q_{j}^{\infty}(\xi+i)^{-j} \tag{B.5.10}
\end{equation*}
$$

Formula (B.5.9) is deduced from identities (B.5.5) - (B.5.7) by Taylor expansion at $\xi+i=$ $\infty$, which allows to expand the functions

$$
\xi^{-j}, \quad\left(\frac{\xi+i 0}{\xi+i}\right)^{\frac{\nu}{2}+\delta} \text { and }\left(\frac{\xi-i 0}{\xi-i}\right)^{\frac{\nu}{2}-\delta}
$$

in negative powers of $(\xi+i)$.
There is also an expansion for

$$
\begin{equation*}
q^{-\infty} \sim \frac{1}{q^{\infty}(\xi)} \sim 1+\sum_{j \geq 1} q_{j}^{-\infty}(\xi+i)^{-j} \tag{B.5.11}
\end{equation*}
$$

so that $q^{\infty} q^{-\infty}$ is a symbol of order $-\infty$.
The following result is well known from Eskin's version [Es1] of the Wiener-Hopf method:
Proposition B.5.4 For $h \in \mathrm{H}^{s-\nu}\left(\mathbb{R}_{+}\right)$, the equation

$$
\begin{equation*}
p_{+} a^{\infty}(D) v=h \tag{B.5.12}
\end{equation*}
$$

has a unique solution $v \in \widetilde{\mathrm{H}}^{s}\left(\mathbb{R}_{+}\right)$. This solution is given by

$$
\begin{equation*}
v=(D+i)^{-\frac{\nu}{2}-\delta} p_{+}(D-i)^{-\frac{\nu}{2}+\delta}\left[a_{0}^{+}\right]^{-1} \tilde{h}, \tag{B.5.13}
\end{equation*}
$$

where $\tilde{h} \in \mathrm{H}^{s-\nu}(\mathbb{R})$ is an extension of $h$ to the whole line.
For $K \in \mathbb{N}$ and $h \in \mathrm{H}^{s-\nu+K}\left(\mathbb{R}_{+}\right)$, this solution $v$ has the asymptotic expansion

$$
\begin{equation*}
v(x)=\sum_{k=0}^{K-1} \chi_{+}(x) x^{\frac{\nu}{2}+\delta+k} e^{-x} d_{k}+v_{\mathrm{rem}, K}(x) \tag{B.5.14}
\end{equation*}
$$

with the remainder $v_{\mathrm{rem}, K} \in \widetilde{\mathrm{H}}^{s+K}\left(\mathbb{R}_{+}\right)$given by

$$
\begin{equation*}
v_{\mathrm{rem}, K}=(D+i)^{-\frac{\nu}{2}-\delta-k} p_{+}(D-i)^{-\frac{\nu}{2}+\delta+k}\left[a_{0}^{+}\right]^{-1} \tilde{h}, \tag{B.5.15}
\end{equation*}
$$

and the coefficients $d_{k}$ by

$$
\begin{equation*}
d_{k}=\frac{e^{-i \frac{\pi}{2}\left(\frac{\nu}{2}+\delta+k+1\right)}}{\Gamma\left(\frac{\nu}{2}+\delta+k+1\right)}\left((D-i)^{-\frac{\nu}{2}+\delta+k}\left(\left[a_{0}^{+}\right]^{-1} \tilde{h}\right)\right)(0) \tag{B.5.16}
\end{equation*}
$$

Let now $u \in \widetilde{\mathrm{H}}^{s}\left(\mathbb{R}_{+}\right)$be a solution of (B.5.4). For $K \in \mathbb{N}$, let $v$ be defined by

$$
\begin{equation*}
v=q^{K}(D) u \quad \text { with } \quad q^{K}(\xi)=1+\sum_{j=1}^{K-1} q_{j}^{\infty}(\xi+i)^{-j} \tag{B.5.17}
\end{equation*}
$$

Note that $(D+i)^{-j}$ is a convolution operator with kernel $\frac{(-i)^{j}}{(j-1)!} x^{j-1} e^{-x} \chi_{+}(x), c f$ Lemma
B.5.1 (i). Thus $v \in \widetilde{\mathrm{H}}^{s}\left(\mathbb{R}_{+}\right)$, and $v$ is solution of

$$
p_{+} a^{\infty}(D) v=g-p_{+} a^{\infty}(D)\left(q^{\infty}(D)-q^{K}(D)\right) u=: h \in \mathrm{H}^{s-\nu+K}\left(\mathbb{R}_{+}\right) .
$$

Therefore $v$ has the expansion (B.5.14), and we can recover the expansion of

$$
\begin{align*}
u & \equiv q^{-\infty}(D)\left(v+\left(q^{\infty}(D)-q^{K}(D)\right) u\right) \quad\left(\bmod \mathscr{C}^{\infty}\right)  \tag{B.5.18}\\
& \equiv q^{-K}(D) v \quad\left(\bmod \widetilde{\mathrm{H}}^{s+K}\left(\mathbb{R}_{+}\right)\right) \quad \text { with } \quad q^{-K}(\xi)=1+\sum_{j=1}^{K-1} q_{j}^{-\infty}(\xi+i)^{-j}
\end{align*}
$$

by simply integrating (B.5.14):

$$
\begin{equation*}
(D+i)^{-j}\left[\chi_{+}(x) x^{\frac{\nu}{2}+\delta+k} e^{-x}\right]=(-1)^{j} \frac{\Gamma\left(\frac{\nu}{2}+\delta+k+1\right)}{\Gamma\left(\frac{\nu}{2}+\delta+k+j+1\right)} \chi_{+}(x) x^{\frac{\nu}{2}+\delta+k+j} e^{-x} \tag{B.5.19}
\end{equation*}
$$

except if $\frac{\nu}{2}+\delta+k \in\{-1,-2, \ldots,-j\}$, where logarithms will appear.
Thus we obtain the asymptotics of $u$ up to regularity $\widetilde{\mathrm{H}}^{s+K}\left(\mathbb{R}_{+}\right)$, and since we assumed that $\frac{\nu}{2}+\delta$ is not a negative integer, no logarithm will appear. We have shown that condition $\left(\mathfrak{H}_{\text {B5 }}\right)$ implies that the asymptotics of $u$ is free of logarithms.
(ii) Let us show the converse. We assume that the equality in $\left(\mathfrak{H}_{\mathrm{B} 5}\right)$ is violated for some $j \geq 1$. Let $M$ be the first such $j$, so that

$$
\begin{equation*}
a(\xi)=a_{0}(\xi) q^{M}(\xi)+a_{M+1}(\xi) \tag{B.5.20}
\end{equation*}
$$

with

$$
q^{M}(\xi)=1+\sum_{j=1}^{M-1} q_{j} \xi^{-j}+\left(q_{M}^{+} \chi_{+}(\xi)+q_{M}^{-} \chi_{-}(\xi)\right)|\xi|^{-M}
$$

and $a_{M+1}(\xi)=\mathscr{O}\left(|\xi|^{-M-1}\right)$ as $|\xi| \rightarrow \infty$.
We will show that there exist $g \in \mathrm{H}^{s-\nu+M+1}\left(\mathbb{R}_{+}\right)$and $u \in \widetilde{\mathrm{H}}^{s}\left(\mathbb{R}_{+}\right)$solution of (B.5.4) such that

$$
\begin{equation*}
u(x)=c_{0} \chi_{+}(x) x^{\frac{\nu}{2}+\delta}+c_{M} \chi_{+}(x) x^{\frac{\nu}{2}+\delta+M} \log x \quad \text { near } x=0 . \tag{B.5.21}
\end{equation*}
$$

The question of regularity of $g=p_{+} a(D) u$ is local at $x=0$. We can therefore stay within the framework of (quasi-)homogeneous distributions and homogeneous symbols, discard lower order terms such as $a_{M+1}(\xi)$, and replace $\xi^{-j}$ by $(\xi+i 0)^{-j}$.
Since the Fourier transform of $\chi_{+}(x) x^{\gamma}$ is $c(\xi+i 0)^{-1-\gamma}$, and the Fourier transform of $\chi_{+}(x) x^{\gamma} \log x$ is $(\xi+i 0)^{-1-\gamma}(c \log (\xi+i 0)+d)$, see Lemma B.5.1, we shall construct the Fourier transform $\hat{u}$ of $u$ in the form

$$
\begin{equation*}
\hat{u}(\xi)=(\xi+i 0)^{-\frac{\nu}{2}-\delta-1}+\hat{c}_{M}(\xi+i 0)^{-\frac{\nu}{2}-\delta-M-1} \log (\xi+i 0)+\hat{d}_{M}(\xi+i 0)^{-\frac{\nu}{2}-\delta-M-1} \tag{B.5.22}
\end{equation*}
$$

We shall show that there exists $\hat{c}_{M} \neq 0$ (hence $c_{M} \neq 0$ ) such that

$$
\begin{equation*}
p_{+} a_{0}(D) q^{M}(D) u \in \mathrm{H}_{\mathrm{loc}}^{s-\nu+M+1}\left(\overline{\mathbb{R}}_{+}\right) . \tag{B.5.23}
\end{equation*}
$$

Since there holds $p_{+}(D-i 0)^{\frac{\nu}{2}-\delta}\left(1-p_{+}\right)=0$, we have the identities

$$
\begin{aligned}
p_{+} a_{0}(D) q^{M}(D) u & =p_{+} a_{0}^{+}(D-i 0)^{\frac{\nu}{2}-\delta}(D+i 0)^{\frac{\nu}{2}+\delta} q^{M}(D) u \\
& =p_{+} a_{0}^{+}(D-i 0)^{\frac{\nu}{2}-\delta} p_{+}(D+i 0)^{\frac{\nu}{2}+\delta} q^{M}(D) u .
\end{aligned}
$$

Therefore, if we prove that $p_{+}(D+i 0)^{\frac{\nu}{2}+\delta} q^{M}(D) u$ belongs to $\mathrm{H}_{\mathrm{loc}}^{s-\frac{\nu}{2}-\delta+M+1}\left(\overline{\mathbb{R}}_{+}\right)$, we have proved (B.5.23).
Consider therefore the Fourier transform $w(\xi)$ of $(D+i 0)^{\frac{\nu}{2}+\delta} q^{M}(D) u$ if $\hat{u}$ has the form (B.5.22):

$$
\begin{aligned}
w(\xi) & =(\xi+i 0)^{\frac{\nu}{2}+\delta} q^{M}(\xi) \hat{u}(\xi) \\
& =\hat{c}_{M}(\xi+i 0)^{-M-1} \log (\xi+i 0)+\left(q_{M}^{+} \chi_{+}(\xi)+q_{M}^{-} \chi_{-}(\xi)\right)|\xi|^{-M}(\xi+i 0)^{-1}+w_{M}(\xi)
\end{aligned}
$$

where
$w_{M}(\xi)=\sum_{j=0}^{M-1} q_{j}(\xi+i 0)^{-j-1}+\hat{d}_{M}(\xi+i 0)^{-M-1}+\hat{c}_{M} \log (\xi+i 0) \mathscr{O}\left(|\xi|^{-M-2}\right)+\mathscr{O}\left(|\xi|^{-M-2}\right)$.
Thus we can discard $w_{M}$, because $p_{+} \mathscr{F}^{-1} w_{M}$ is sufficiently regular.
Now we use the additive decomposition, see Lemma B.5.2, for $\xi \neq 0$ :

$$
\begin{align*}
& \left(q_{M}^{+} \chi_{+}(\xi)+q_{M}^{-} \chi_{-}(\xi)\right)|\xi|^{-M}(\xi+i 0)^{-1}=q_{M}^{+}(\xi+i 0)^{-M-1}+  \tag{B.5.25}\\
& \quad \frac{1}{2 i \pi}\left((-1)^{M} q_{M}^{-}-q_{M}^{+}\right)\left((\xi+i 0)^{-M-1} \log (\xi+i 0)-(\xi-i 0)^{-M-1} \log (\xi-i 0)\right)
\end{align*}
$$

The only non-regular contribution to $p_{+} \mathscr{F}^{-1} w$ comes from the term

$$
(\xi+i 0)^{-M-1} \log (\xi+i 0)\left[\hat{c}_{M}+\frac{1}{2 i \pi}\left((-1)^{M} q_{M}^{-}-q_{M}^{+}\right)\right]
$$

This term is absent if

$$
\begin{equation*}
\hat{c}_{M}+\frac{1}{2 i \pi}\left((-1)^{M} q_{M}^{-}-q_{M}^{+}\right)=0 \tag{B.5.26}
\end{equation*}
$$

We see that the possibility of having $\hat{c}_{M} \neq 0$ together with condition (B.5.26) is a consequence of the violation of equality $\left(\mathfrak{H}_{\mathrm{B} 5}\right)$ for $j=M$. The proof is complete.

## B. 6 AUXILIARY RESULTS ON $\Psi$ DO

We need some results for pseudodifferential operators ( $\Psi D O$ ) of one variable acting on functions of $n$ variables, and also the connection between $\Psi D O$ in $n$ variables and reduced $\Psi D O$ in one variable. The suitable function spaces were introduced in Section B.1. Here, we only need the "model" domain for the boundary of $\mathscr{M}$, that is $\mathbb{R}_{+}^{n}=\mathbb{R}^{n-1} \times \mathbb{R}_{+}$with coordinates $x=\left(x^{\prime}, x_{n}\right)$ and dual coordinates $\xi=\left(\xi^{\prime}, \xi_{n}\right)$.

The following lemma is a particular case of Theorem 1.11 and Lemma 2.9 in [ChkDu1].

Lemma B.6.1 Let the symbol $\mathrm{b}=\mathrm{b}\left(x ; \xi_{n}\right)$ satisfy $\partial_{x}^{\alpha} \partial_{\xi_{n}}^{k} \mathrm{~b}\left(x ; \xi_{n}\right)=\mathscr{O}\left(\left|\xi_{n}\right|^{\nu-k}\right)$ as $\left|\xi_{n}\right| \rightarrow$ $\infty$ for all $k \in \mathbb{N}_{0}, \alpha \in \mathbb{N}_{0}^{n}, x \in \mathbb{R}^{n}, \xi_{n} \in \mathbb{R}$. Let $\kappa \in \mathbb{N}_{0}, s \in \mathbb{R}$. Then the pseudodifferential operator $\mathrm{b}\left(x ; D_{n}\right)$ is bounded between anisotropic Bessel potential spaces:

$$
\begin{equation*}
\mathrm{b}\left(x ; D_{n}\right): \mathrm{H}^{(\infty, s), \kappa}\left(\mathbb{R}^{n}\right) \longrightarrow \mathrm{H}^{(\infty, s-\nu), \kappa}\left(\mathbb{R}^{n}\right) . \tag{B.6.1}
\end{equation*}
$$

If, in addition, $\operatorname{supp} \mathrm{b}(x, \cdot)$ is compact, for all $x \in \mathbb{R}^{n}$, then the operator $\mathrm{b}\left(x ; D_{n}\right)$ is a smoothing operator:

$$
\begin{equation*}
\mathrm{b}\left(x ; D_{n}\right): \mathrm{H}^{(\infty, s), \kappa}\left(\mathbb{R}^{n}\right) \longrightarrow \mathscr{C}^{\infty}\left(\mathbb{R}^{n}\right) \tag{B.6.2}
\end{equation*}
$$

Proof. It is easy to check that the operator

$$
\begin{equation*}
\mathrm{b}\left(x ; D_{n}\right): \mathrm{H}^{(\mu, s), \kappa}\left(\mathbb{R}^{n}\right) \longrightarrow \mathrm{H}^{(\mu-\sigma(\nu), s-\nu), \kappa}\left(\mathbb{R}^{n}\right) . \tag{B.6.3}
\end{equation*}
$$

is bounded, where

$$
\sigma(\nu):= \begin{cases}0 & \text { for } \quad \nu>0 \\ |\nu| & \text { for } \quad \nu<0\end{cases}
$$

In fact, the boundedness (B.6.3) follows from the Mikhlin-Hörmander theorem on multipliers since

$$
\forall|\alpha| \leq n, \quad \forall \xi=\left(\xi^{\prime}, \xi_{n}\right) \in \mathbb{R}^{n}, \quad \xi^{\alpha} \partial_{\xi}^{\alpha}\left[\frac{\left\langle\xi^{\prime}\right\rangle^{\mu-\sigma(\nu)}\langle\xi\rangle^{s-\nu}\langle\xi\rangle^{\nu}}{\left\langle\xi^{\prime}\right\rangle^{\mu}\langle\xi\rangle^{s}}\right] \leq 1 .
$$

The boundedness (B.6.1) is a consequence of (B.6.3).
As for (B.6.2), it follows from (B.6.1) because the symbol b satisfies $\partial_{x}^{\alpha} \partial_{\xi_{n}}^{k} \mathrm{~b}\left(x ; \xi_{n}\right)=$ $\mathcal{O}\left(\left|\xi_{n}\right|^{\nu-k}\right)$ for arbitrary $\nu<0$.

The following lemma generalizes Eskin's Wiener-Hopf technique from the scalar onedimensional case, see Proposition B.5.4, to systems of multidimensional pseudodifferential equations.

Lemma B.6.2 Let us consider the principal part $\mathbf{a}_{\mathrm{pr}}$ of the symbol $\mathbf{a}$ in (B.3.1) with the ellipticity condition ( $\mathfrak{H}_{\mathrm{B} 1}$ ). We introduce

$$
\begin{equation*}
\mathbf{a}^{\infty}\left(x^{\prime} ; \xi_{n}\right):=\left\langle\xi_{n}\right\rangle^{\nu} \mathbf{a}_{\mathrm{pr}}\left(x^{\prime}, 0 ; 0,+1\right) . \tag{B.6.4}
\end{equation*}
$$

Let s, $\nu \in \mathbb{R}$ such that $\frac{\nu}{2}-\frac{1}{2}<s<\frac{\nu}{2}+\frac{1}{2}, \kappa \in \mathbb{N}_{0}$. Then the system of equations

$$
\begin{equation*}
\boldsymbol{p}_{+} \mathbf{a}^{\infty}\left(x^{\prime} ; D_{n}\right) \boldsymbol{u}=\boldsymbol{g}, \quad \boldsymbol{g} \in \mathrm{H}^{(\infty, s-\nu), \kappa}\left(\mathbb{R}_{+}^{n}\right)^{N} \tag{B.6.5}
\end{equation*}
$$

where $\boldsymbol{p}_{+}$is the restriction from $\mathbb{R}^{n}$ to $\mathbb{R}_{+}^{n}$, has a unique solution $\boldsymbol{u} \in \widetilde{\mathrm{H}}^{(\infty, s), \kappa}\left(\mathbb{R}_{+}^{n}\right)^{N}$, represented by the formula

$$
\begin{equation*}
\boldsymbol{u}=\left(D_{n}+i\right)^{-\frac{\nu}{2}} \chi_{+}\left(D_{n}-i\right)^{-\frac{\nu}{2}}\left[\mathbf{a}_{\mathrm{pr}}\left(x^{\prime}, 0 ; 0,1\right)\right]^{-1} \boldsymbol{g} \tag{B.6.6}
\end{equation*}
$$

where $\chi_{+}\left(x_{n}\right)$ is the characteristic function of the half space $\mathbb{R}_{+}^{n} \subset \mathbb{R}^{n}$.

For arbitrary $K \in \mathbb{N}, K \leq \kappa$, and $\boldsymbol{g} \in \mathrm{H}^{(\infty, s-\nu+K), \kappa}\left(\mathbb{R}_{+}^{n}\right)$ this solution has the following asymptotic expansions

$$
\begin{aligned}
\boldsymbol{u}\left(x^{\prime}, x_{n}\right) & =\sum_{k=0}^{K-1} x_{n}^{\frac{\nu}{2}+k} e^{-x_{n}} \boldsymbol{d}^{k}\left(x^{\prime}\right)+\boldsymbol{u}_{\mathrm{rem}, K}\left(x^{\prime}, x_{n}\right), \quad \boldsymbol{u}_{\mathrm{rem}, K} \in \widetilde{\mathrm{H}}^{(\infty, s+K), \kappa}\left(\mathbb{R}_{+}^{n}\right)^{N} \\
& =\sum_{k=0}^{K-1} x_{n}^{\frac{\nu}{2}+k} \boldsymbol{d}_{0}^{k}\left(x^{\prime}\right)+\boldsymbol{u}_{\mathrm{rem}, K}^{0}\left(x^{\prime}, x_{n}\right), \quad \boldsymbol{u}_{\mathrm{rem}, K}^{0} \in \widetilde{\mathrm{H}}^{(\infty, s+K), \kappa}\left(\mathbb{R}_{+}^{n}\right)^{N} .
\end{aligned}
$$

with the $\mathscr{C}^{\infty}\left(\mathbb{R}^{n-1}\right)$ coefficients

$$
\begin{equation*}
\boldsymbol{d}^{k}\left(x^{\prime}\right):=\frac{e^{-\frac{\pi}{2} i\left(\frac{\nu}{2}+k+1\right)}}{\Gamma\left(\frac{\nu}{2}+k+1\right)}\left(\left(D_{n}-i\right)^{-\frac{\nu}{2}+k}\left[\mathbf{a}_{\mathrm{pr}}\left(\left(x^{\prime}, 0 ; 0,1\right)\right]^{-1} \boldsymbol{g}\right)\left(x^{\prime}, 0\right),\right. \tag{B.6.7}
\end{equation*}
$$

and

$$
\boldsymbol{d}_{0}^{k}\left(x^{\prime}\right):=\sum_{\ell=0}^{k} \frac{(-1)^{k-\ell}}{(k-\ell)!} \boldsymbol{d}^{\ell}\left(x^{\prime}\right), \quad k=1, \ldots, M
$$

For the proof see [ChkDu1, Lemma 2.6]. Note that, by its mere definition, $\mathbf{a}^{\infty}$ satisfies itself condition $\left(\mathfrak{H}_{\mathrm{B} 4}\right)$.

The following Lemma B. 6.3 will serve for the evaluation of the terms and the remainders in the Taylor expansions which will provide the next Lemma B.6.4.

Lemma B.6.3 Let $\mathrm{b}(x ; D)$ be a $\Psi D O$ such that for an integer $\bar{m} \in \mathbb{N}$ its symbol satisfies $\mathrm{b}(x ; \xi)=x_{n}^{\bar{m}} \mathrm{~b}(x ; \xi)$ with b in the class $S_{\mathrm{cl}}^{\nu}\left(\mathbb{R}_{+}^{n} \times \mathbb{R}^{n}\right)$. We suppose that, moreover, there exists an integer $\bar{k} \in \mathbb{N}_{0}$ such that $\partial_{x}^{\alpha} \partial_{\xi}^{\gamma} \breve{\mathrm{b}}(x ; \xi)=\mathscr{O}\left(\left|\xi^{\prime}\right|^{\bar{k}-\left|\gamma^{\prime}\right|}\left|\xi_{n}\right|^{\nu-\bar{k}-\gamma_{n}}\right)$ for all $\alpha$ and $\gamma=\left(\gamma^{\prime}, \gamma_{n}\right) \in \mathbb{N}^{n}$. Then for all $\mu, s \in \mathbb{R}$ and $\kappa \geq m, \mathrm{~b}(x ; D)$ is bounded between the spaces:

$$
\begin{equation*}
\mathrm{b}(x ; D): \mathrm{H}^{(\mu, s), \kappa}\left(\mathbb{R}^{n}\right) \longrightarrow \mathrm{H}^{(\mu-\bar{k}, s+\bar{k}+\bar{m}-\nu), \kappa-\bar{m}}\left(\mathbb{R}^{n}\right) \tag{B.6.8}
\end{equation*}
$$

Lemma B.6.4 Let $j \in \mathbb{N}_{0}$ and let us consider the homogeneous part $\mathbf{a}_{j}$ of degree $\nu-j$ of the symbol $\mathbf{a}$ in (B.3.1). For any $K \in \mathbb{N}$, there holds the expansion of the symbol $\mathbf{a}_{j}$

$$
\begin{equation*}
\mathbf{a}_{j}(x ; \xi)=\sum_{m+\left|\gamma^{\prime}\right| \leq K-1} x_{n}^{m}\left(\xi^{\prime}\right)^{\gamma^{\prime}} \breve{\mathbf{a}}_{j ; m, \gamma^{\prime}}\left(x^{\prime} ; \operatorname{sgn} \xi_{n}\right) \xi_{n}^{-j-\left|\gamma^{\prime}\right|}\left|\xi_{n}\right|^{\nu}+\mathbf{a}_{j ; \mathrm{rem}, K} \tag{B.6.9}
\end{equation*}
$$

with $\breve{\mathbf{a}}_{j ; m, \gamma^{\prime}}\left(x^{\prime} ; \omega\right)=\frac{1}{m!} \frac{1}{\gamma^{\prime}!} \omega^{j+\left|\gamma^{\prime}\right|} \partial_{x_{n}}^{m} \partial_{\xi^{\prime}}^{\gamma^{\prime}} \mathbf{a}_{j}\left(x^{\prime}, 0 ; 0, \omega\right), x^{\prime} \in \mathbb{R}^{n-1}, \omega= \pm 1$, and $\mathbf{a}_{j ; \mathrm{rem}, K}$ bounded between the spaces

$$
\begin{equation*}
\mathbf{a}_{j ; \mathrm{rem}, K}(x ; D): \mathrm{H}^{(\infty, s), \infty}\left(\mathbb{R}^{n}\right)^{N} \longrightarrow \mathrm{H}^{(\infty, s+K-\nu), \infty}\left(\mathbb{R}^{n}\right)^{N} \tag{B.6.10}
\end{equation*}
$$

If condition $\left(\mathfrak{H}_{\mathrm{B} 4}\right)$ holds, then $\breve{\mathbf{a}}_{j ; m, \gamma^{\prime}}\left(x^{\prime} ; \omega\right)=\breve{\mathbf{a}}_{j ; m, \gamma^{\prime}}\left(x^{\prime}\right)$ does not depend on $\omega$.

Proof. The Taylor formula, applied at $x_{n}=0$, and then at $\left|\xi_{n}\right|^{-1} \xi^{\prime}=0$, gives:

$$
\begin{array}{r}
\mathbf{a}_{j}\left(x^{\prime}, r ; \xi^{\prime}, \xi_{n}\right)= \\
\sum_{m=0}^{K-1} \frac{x_{n}^{m}}{m!}\left(\partial_{x_{n}}^{m} \mathbf{a}_{j}\right)\left(x^{\prime}, 0 ;\left|\xi_{n}\right|^{-1} \xi^{\prime}, \operatorname{sgn} \xi_{n}\right)\left|\xi_{n}\right|^{\nu-j}+x_{n}^{K} \mathbf{a}_{j ; \text { rem }, K}^{(1)}(x ; \xi) \\
=\sum_{m=0}^{K-1} \frac{x_{n}^{m}}{m!} \sum_{\left|\gamma^{\prime}\right|=0}^{K-1-m}\left|\xi_{n}\right|^{\nu-j-\left|\gamma^{\prime}\right|} \frac{\left(\xi^{\prime}\right)^{\gamma^{\prime}}}{\left(\gamma^{\prime}\right)!}\left(\partial_{x_{n}}^{m} \partial_{\xi^{\prime}}^{\gamma^{\prime}} \mathbf{a}_{j}\right)\left(x^{\prime}, 0 ; 0, \operatorname{sgn} \xi_{n}\right) \\
+\mathbf{a}_{j ; \text { rem }, K}(x ; \xi),
\end{array}
$$

where the remainder can be written as

$$
\mathbf{a}_{j ; \mathrm{rem}, K}(x ; \xi)=x^{K} \mathbf{a}_{j ; \mathrm{rem}, K}^{(1)}(x ; \xi)+\sum_{m=0}^{K-1} x_{n}^{m} \mathbf{a}_{j ; \mathrm{rem}, m, K-m}^{(2)}\left(x^{\prime} ; \xi\right)
$$

where $x_{n}^{K} \mathbf{a}_{j ; \text { rem, } K}^{(1)}$ satisfies the assumptions of Lemma B.6.3 with $\bar{m}=K$ and $\bar{k}=0$, and $x_{n}^{m} \mathbf{a}_{j ; \text { rem }, m, K-m}^{(2)}$ with $\bar{m}=m$ and $\bar{k}=K-m$. Taking $\mu=\infty$ and $\kappa=\infty$, we obtain the lemma.

A standard Taylor expansion of the function $\left\langle\xi_{n}\right\rangle^{\nu}$ at $\xi_{n}= \pm \infty$ yields the following expansion of the symbol $\mathbf{a}^{\infty}$ (B.6.4):

Lemma B.6.5 Let us consider the symbol $\mathbf{a}^{\infty}$ defined in (B.6.4). For any integer $K \in \mathbb{N}$, there holds the expansion

$$
\begin{equation*}
\mathbf{a}^{\infty}\left(x^{\prime} ; \xi_{n}\right)=\sum_{j \leq K-1} \breve{\mathbf{a}}_{j}^{\infty}\left(x^{\prime}\right) \xi_{n}^{-j}\left|\xi_{n}\right|^{\nu}+\mathbf{a}_{\mathrm{rem}, K}^{\infty}\left(x^{\prime} ; \xi_{n}\right) \tag{B.6.11}
\end{equation*}
$$

with $\breve{\mathbf{a}}_{0}^{\infty}\left(x^{\prime}\right)=\mathbf{a}_{0}\left(x^{\prime}, 0 ; 0,+1\right), \breve{\mathbf{a}}_{j}^{\infty}\left(x^{\prime}\right)=c_{j} \mathbf{a}_{0}\left(x^{\prime}, 0 ; 0,+1\right)$ with $c_{j} \in \mathbb{R}$, and $\mathbf{a}_{\mathrm{rem}, K}^{\infty}$ is a bounded operator between the spaces

$$
\begin{equation*}
\mathrm{a}_{\mathrm{rem}, K}^{\infty}\left(x^{\prime} ; D_{n}\right): \mathrm{H}^{(\infty, s), \infty}\left(\mathbb{R}^{n}\right)^{N} \longrightarrow \mathrm{H}^{(\infty, s+K-\nu), \infty}\left(\mathbb{R}^{n}\right)^{N} \tag{B.6.12}
\end{equation*}
$$

## B. 7 Proof of the main theorem of Part B

We are going to prove Theorem B.4.1. Let us start by reformulation of the conditions of equation (B.3.2): we consider

$$
\begin{equation*}
\phi \in \widetilde{\mathrm{H}}^{(\mu, s), \kappa}(\mathscr{M})^{N} \quad \text { such that } \quad \mathbf{a}\left(\mathscr{X} ; D_{\mathscr{X}}\right) \boldsymbol{\phi}=\boldsymbol{g}, \quad \text { with } \quad \boldsymbol{g} \in \mathrm{H}^{(\mu, s-\nu), \kappa}(\mathscr{M})^{N} \tag{B.7.1}
\end{equation*}
$$

for arbitrary $-\infty<\mu \leq \infty$. In [ChkDu1, Theorem 1.12] it is proved that the system (B.3.2) is Fredholm (or is uniquely solvable) if and only if the system (B.7.1) is Fredholm (is uniquely solvable) and these equations have equal dimensions of kernels and cokernels.

Since the assertion is local, we can suppose that our domain is the half-space $\mathbb{R}_{+}^{n}$, and all functions and symbols are compactly supported in the variable $x \in \mathbb{R}_{+}^{n}$. We recall that $x=\left(x^{\prime}, x_{n}\right)$ and its dual variable is $\xi=\left(\xi^{\prime}, \xi_{n}\right)$.

Homogeneous symbols and the kernels of the corresponding $\Psi D O$ with negative order have singularities at 0 . Multiplying them by a function $\chi^{0} \in \mathscr{C}^{\infty}(\mathbb{R})$, where $\chi^{0}\left(\xi_{n}\right)=0$ for $\left|\xi_{n}\right|<1$ and $\chi^{0}\left(\xi_{n}\right)=1$ for $\left|\xi_{n}\right|>2$ we cut the singularity off. The perturbation operator is smoothing: $\left[I-\chi^{0}\left(D_{n}\right)\right] \varphi \in \mathscr{C}_{0}^{\infty}\left(\mathbb{R}^{n}\right)$ for arbitrary $\varphi \in \mathrm{H}^{(\infty, \mu), \kappa}\left(\mathbb{R}^{n}\right)$ (see Lemma (B.6.1)), and will be ignored. Although we do not write the cut off function, we suppose it is present and forget about singularities of symbols at $\xi_{n}=0$.

Since $\mathscr{C}_{0}^{\infty}\left(\overline{\mathbb{R}_{+}^{n}}\right)^{N} \subset \mathrm{H}^{(\infty, s-\nu+M+1), \infty}\left(\mathbb{R}_{+}^{n}\right)^{N}$, for any $M \in \mathbb{N}_{0}$, it is sufficient to derive the asymptotics for a solution of equation (B.7.1). Relying on the expansion of the classical symbol $\mathbf{a}(x ; \xi)$ :

$$
\begin{gather*}
\mathbf{a}=\sum_{j=0}^{M} \mathbf{a}_{j}+\mathbf{a}_{\mathrm{rem}, M+1},  \tag{B.7.2}\\
\mathbf{a}_{j} \in \mathrm{~S}_{\mathrm{hom}}^{\nu-j}\left(\mathbb{R}_{+}^{n} \times \mathbb{R}^{n}\right)^{N \times N}, \quad \mathbf{a}_{\mathrm{rem}, M+1} \in \mathrm{~S}_{\mathrm{cl}}^{\nu-M-1}\left(\mathbb{R}_{+}^{n} \times \mathbb{R}^{n}\right)^{N \times N}
\end{gather*}
$$

we will apply induction on $M$, starting with the case $M=0$.
For $M=0$, the equation (B.7.1) (with $\mathscr{M}=\mathbb{R}_{+}^{n}$, as agreed) is written in the following equivalent form

$$
\begin{equation*}
\boldsymbol{p}_{+} \mathbf{a}^{\infty}\left(x^{\prime} ; D_{n}\right) \boldsymbol{\phi}=\boldsymbol{g}^{\infty} \tag{B.7.3}
\end{equation*}
$$

where $\mathbf{a}^{\infty}\left(x^{\prime} ; \xi_{n}\right)$ is defined in (B.6.4) and

$$
\boldsymbol{g}^{\infty}:=\boldsymbol{g}-\mathbf{a}_{\mathrm{rem}, 1}\left(x ; D_{x}\right) \boldsymbol{\phi}-\left[\mathbf{a}_{0}\left(x ; D_{x}\right)-\mathbf{a}^{\infty}\left(x^{\prime} ; D_{n}\right)\right] \boldsymbol{\phi}
$$

We observe
(i) The remainder $\mathbf{a}_{\mathrm{rem}, 1}\left(x ; D_{x}\right): \mathrm{H}^{(\infty, s), \infty}\left(\mathbb{R}^{n}\right)^{N} \longrightarrow \mathrm{H}^{(\infty, s+1-\nu), \infty}\left(\mathbb{R}^{n}\right)^{N}$ is bounded.
(ii) Lemma B.6.4 for $j=0, K=1$ gives that

$$
\mathbf{a}_{0}(x ; \xi)=\mathbf{a}_{0}\left(x^{\prime}, 0 ; 0, \operatorname{sgn} \xi_{n}\right)\left|\xi_{n}\right|^{\nu}+\mathbf{a}_{0 ; \text { rem }, 1}(x, \xi)
$$

with $\mathbf{a}_{0 ; \text { rem }, 1}\left(x ; D_{x}\right): \mathrm{H}^{(\infty, s), \infty}\left(\mathbb{R}^{n}\right)^{N} \longrightarrow \mathrm{H}^{(\infty, s+1-\nu), \infty}\left(\mathbb{R}^{n}\right)^{N}$ bounded.
(iii) Lemma B.6.5 for $j=0, K=1$ gives that

$$
\mathbf{a}^{\infty}\left(x ; \xi_{n}\right)=\mathbf{a}_{0}\left(x^{\prime}, 0 ; 0,+1\right)\left|\xi_{n}\right|^{\nu}+\mathbf{a}_{\text {rem }, 1}^{\infty}\left(x, \xi_{n}\right)
$$

with $\mathbf{a}_{\mathrm{rem}, 1}^{\infty}\left(x ; D_{n}\right): \mathrm{H}^{(\infty, s), \infty}\left(\mathbb{R}^{n}\right)^{N} \longrightarrow \mathrm{H}^{(\infty, s+1-\nu), \infty}\left(\mathbb{R}^{n}\right)^{N}$ bounded.
Condition $\left(\mathfrak{H}_{\mathrm{B} 4}\right)$ yields that $\mathbf{a}_{0}\left(x^{\prime}, 0 ; 0, \operatorname{sgn} \xi_{n}\right)=\mathbf{a}_{0}\left(x^{\prime}, 0 ; 0,+1\right)$, therefore $\mathbf{a}_{0}-\mathbf{a}^{\infty}=$ $\mathbf{a}_{0 ; \text { rem, } 1}-\mathbf{a}_{\text {rem, } 1}^{\infty}$. We deduce from $\phi \in \mathrm{H}^{(\infty, s), \infty}\left(\mathbb{R}^{n}\right)^{N}$, that

$$
\begin{equation*}
\boldsymbol{g}^{\infty} \in \mathrm{H}^{(\infty, s+1-\nu), \infty}\left(\mathbb{R}^{n}\right)^{N} \tag{B.7.4}
\end{equation*}
$$

Invoking Lemma B.6.2 we derive the expansion (B.4.3) for $K=1$ :

$$
\begin{align*}
& \boldsymbol{\phi}=\boldsymbol{\phi}_{0}+\boldsymbol{\phi}_{\mathrm{rem}, 1}, \quad \boldsymbol{\phi}_{0}\left(x^{\prime}, x_{n}\right)=\boldsymbol{d}^{0}\left(x^{\prime}\right) x_{n}^{\frac{\nu}{2}} e^{-x_{n}}  \tag{B.7.5}\\
& \boldsymbol{d}^{0} \in \mathscr{C}^{\infty}\left(\mathbb{R}^{n-1}\right)^{N}, \quad \boldsymbol{\phi}_{\mathrm{rem}, 1} \in \widetilde{\mathrm{H}}^{(\infty, s+1), \infty}\left(\mathbb{R}_{+}^{n}\right)^{N}
\end{align*}
$$

Now let $M \geq 1$ and suppose we have proved

$$
\begin{align*}
& \boldsymbol{\phi}=\sum_{k=0}^{M-1} \boldsymbol{\phi}_{k}+\boldsymbol{\phi}_{\mathrm{rem}, M}, \quad \boldsymbol{\phi}_{k}\left(x^{\prime}, x_{n}\right):=\boldsymbol{d}^{k}\left(x^{\prime}\right) x_{n}^{\frac{\nu}{n}+k} e^{-x_{n}},  \tag{B.7.6}\\
& \boldsymbol{d}^{k} \in \mathscr{C}^{\infty}\left(\mathbb{R}^{n-1}\right)^{N}, \quad \phi_{\mathrm{rem}, M} \in \widetilde{\mathrm{H}}^{(\infty, s+M), \infty}\left(\mathbb{R}_{+}^{n}\right)^{N}
\end{align*}
$$

It can be proved that $\phi_{k} \in \widetilde{\mathrm{H}}^{(\infty, s+k), \infty}\left(\mathbb{R}_{+}^{n}\right)^{N}$, because $\frac{\nu}{2}-s>-\frac{1}{2}$, see [ChkDu1, (2.30)]. Then the right hand side $\boldsymbol{g}^{\infty}$ of equation (B.7.3) can be represented as follows

$$
\begin{align*}
\boldsymbol{g}^{\infty}=\boldsymbol{g}_{\mathrm{rem}, M+1}^{1} & -\sum_{j=1}^{M} \sum_{k=0}^{M-j} \boldsymbol{p}_{+} \mathbf{a}_{j}\left(x ; D_{x}\right) \boldsymbol{\phi}_{k} \\
& -\sum_{k=0}^{M-1} \boldsymbol{p}_{+}\left[\mathbf{a}_{0}\left(x ; D_{x}\right)-\mathbf{a}^{\infty}\left(x^{\prime} ; D_{n}\right)\right] \boldsymbol{\phi}_{k} \tag{B.7.7}
\end{align*}
$$

where

$$
\begin{aligned}
\boldsymbol{g}_{\mathrm{rem}, M+1}^{1}=\boldsymbol{g} & -\boldsymbol{p}_{+} \mathbf{a}_{\mathrm{rem}, M+1}\left(x ; D_{x}\right) \boldsymbol{\phi} \\
& -\sum_{j=1}^{M} \boldsymbol{p}_{+} \mathbf{a}_{j}\left(x ; D_{x}\right) \boldsymbol{\phi}_{\mathrm{rem}, M-j+1} \\
& -\boldsymbol{p}_{+}\left[\mathbf{a}_{0}\left(x ; D_{x}\right)-\mathbf{a}^{\infty}\left(x^{\prime} ; D_{n}\right)\right] \boldsymbol{\phi}_{\mathrm{rem}, M}
\end{aligned}
$$

It is clear from the arguments used for the step $M=0$ that $\boldsymbol{g}_{\mathrm{rem}, M+1}^{1}$ belongs to the space $\mathrm{H}^{(\infty, s-\nu+M+1), \infty}\left(\mathbb{R}_{+}^{n}\right)^{N}$.

We now use the expansion (B.6.9) with $K=M+1-j-k$ for the term $\mathbf{a}_{j}\left(x ; D_{x}\right) \boldsymbol{\phi}_{k}$, and the expansion (B.6.11) with $K=M+1-k$ for the term $\mathbf{a}^{\infty}\left(x^{\prime} ; D_{n}\right) \phi_{k}$. Taking into account that condition $\left(\mathfrak{H}_{\mathrm{B} 4}\right)$ holds, we obtain

$$
\begin{equation*}
\boldsymbol{g}^{\infty}=\boldsymbol{g}_{\mathrm{rem}, M+1}^{2}-\sum_{k=0}^{M-1} \sum_{\substack{j, m, \gamma^{\prime} \\ 0<j+m+\left|\gamma^{\prime}\right| \leq M-k}} \mathbf{b}_{j ; m, \gamma^{\prime}}\left(x ; D_{x}\right) \boldsymbol{\phi}_{k} \tag{B.7.8}
\end{equation*}
$$

where $\boldsymbol{g}_{\text {rem }, M+1}^{2}$ belongs to $\mathrm{H}^{(\infty, s-\nu+M+1), \infty}\left(\mathbb{R}_{+}^{n}\right)^{N}$ and

$$
\mathbf{b}_{j ; m, \gamma^{\prime}}(x ; \xi)=x_{n}^{m}\left(\xi^{\prime}\right)^{\gamma^{\prime}} \breve{\mathbf{b}}_{j ; m, \gamma^{\prime}}\left(x^{\prime}\right) \xi_{n}^{-j-\left|\gamma^{\prime}\right|}\left|\xi_{n}\right|^{\nu}
$$

with $\breve{\mathbf{b}}_{j ; m, \gamma^{\prime}}\left(x^{\prime}\right)$ defined for $x^{\prime} \in \mathbb{R}^{n-1}$ as follows

$$
\breve{\mathbf{b}}_{j ; m, \gamma^{\prime}}\left(x^{\prime}\right):= \begin{cases}\breve{\mathbf{a}}_{j ; m, \gamma^{\prime}}\left(x^{\prime}\right) & \text { if } m+\left|\gamma^{\prime}\right| \neq 0 \\ \breve{\mathbf{a}}_{j ; 0,0}\left(x^{\prime}\right)-\breve{\mathbf{a}}_{j}^{\infty}\left(x^{\prime}\right) & \text { if } m=0, \gamma^{\prime}=0\end{cases}
$$

Now we use formula (B.6.6) Lemma B.6.2 to invert the operator $\mathbf{a}^{\infty}\left(x^{\prime} ; D_{n}\right)$, and from equation (B.7.3) $\boldsymbol{p}_{+} \mathbf{a}^{\infty}\left(x^{\prime} ; D_{n}\right) \boldsymbol{\phi}=\boldsymbol{g}^{\infty}$ with the expansion (B.7.8) of $\boldsymbol{g}^{\infty}$, we find

$$
\begin{align*}
\boldsymbol{\phi}= & \sum_{k=0}^{M-1} \sum_{\substack{j, m, \gamma^{\prime} \\
0<j+m+\left|\gamma^{\prime}\right| \leq M-k}}\left(D_{n}+i\right)^{-\frac{\nu}{2}} \boldsymbol{p}_{+}\left(D_{n}-i\right)^{-\frac{\nu}{2}} \breve{\mathbf{a}}_{0}^{-1}\left(x^{\prime}\right) \mathbf{b}_{j ; m, \gamma^{\prime}}\left(x ; D_{x}\right) \boldsymbol{\phi}_{k}  \tag{B.7.9}\\
& +\left[\boldsymbol{p}_{+}\left(\mathbf{a}^{\infty}\left(x^{\prime} ; D_{n}\right)\right]^{-1} \boldsymbol{g}_{\mathrm{rem}, M+1}^{2},\right.
\end{align*}
$$

where $\breve{\mathbf{a}}_{0}\left(x^{\prime}\right)=\mathbf{a}_{\mathrm{pr}}\left(x^{\prime}, 0 ; 0,1\right)$. Recalling Lemma B.5.1 (i), and using a Taylor expansion at $\xi_{n}=\infty$, we find the following, $c f$ (B.7.6) for the $\phi_{k}$ :

$$
\begin{align*}
& \mathscr{F}_{x_{n} \rightarrow \xi_{n}}\left[\boldsymbol{\phi}_{k}\left(x^{\prime}, x_{n}\right)\right]=\mathscr{F}_{x_{n} \rightarrow \xi_{n}}\left[x_{n}^{\frac{\nu}{2}+k} \boldsymbol{d}^{k}\left(x^{\prime}\right) e^{-x_{n}}\right] \\
& \quad=\left(\xi_{n}+i\right)^{-\frac{\nu}{2}-k-1} e^{\frac{\pi}{2}\left(\frac{\nu}{2}+k+1\right) i} \Gamma\left(\frac{\nu}{2}+k+1\right) \boldsymbol{d}^{k}\left(x^{\prime}\right) \\
& \quad=\sum_{q=0}^{M} \boldsymbol{d}^{k q}\left(x^{\prime}\right)\left(\xi_{n}+i 0\right)^{-\frac{\nu}{2}-k-q-1}+\left(\xi_{n}+i\right)^{-\frac{\nu}{2}-M-q-2} \boldsymbol{d}_{\text {rem }, M}^{k}\left(x^{\prime}\right) \tag{B.7.10}
\end{align*}
$$

and the last summand is ignored in the sequel because it contributes into the smooth remainder term. From (B.7.9) and (B.7.10), we see that modulo a remainder $\phi_{\text {rem; } M+1}^{1}$ in the space $\mathrm{H}^{(\infty, s+M+1), \infty}\left(\mathbb{R}_{+}^{n}\right)^{N}, \phi$ is a finite sum of terms $\varphi$ which have the generic form

$$
\begin{align*}
\boldsymbol{\varphi} & =\left(D_{n}+i\right)^{-\frac{\nu}{2}} \boldsymbol{\psi} \quad \text { with }  \tag{B.7.11}\\
\boldsymbol{\psi} & =\boldsymbol{p}_{+}\left(D_{n}-i\right)^{-\frac{\nu}{2}} \breve{a}_{0}^{-1}\left(x^{\prime}\right) \boldsymbol{h}(x) \quad \text { with }  \tag{B.7.12}\\
\boldsymbol{h} & =x_{n}^{m} \mathscr{F}_{\xi \rightarrow x}\left\{\left(\xi^{\prime}\right)^{\gamma^{\prime}} \breve{\mathbf{b}}\left(x^{\prime}\right) \xi_{n}^{-\ell}\left|\xi_{n}\right|^{\nu} \times \mathscr{F}_{x^{\prime} \rightarrow \xi^{\prime}}\left[\boldsymbol{d}\left(x^{\prime}\right)\right]\left(\xi_{n}+i 0\right)^{-\frac{\nu}{2}-q-1}\right\} \tag{B.7.13}
\end{align*}
$$

for $m, \ell, q \in \mathbb{N}_{0}, \gamma^{\prime} \in \mathbb{N}^{n-1}$, and $\breve{\mathbf{b}} \in \mathscr{C}^{\infty}\left(\mathbb{R}^{n-1}\right)^{N \times N}, \boldsymbol{d} \in \mathscr{C}^{\infty}\left(\mathbb{R}^{n-1}\right)$. Let us study $h$ first:

$$
\begin{align*}
\boldsymbol{h}(x) & =x_{n}^{m} \breve{\mathbf{b}}\left(x^{\prime}\right)\left[\left(i \partial_{x^{\prime}}\right)^{\gamma^{\prime}} \boldsymbol{d}\right]\left(x^{\prime}\right) \times \mathscr{F}_{\xi_{n} \rightarrow x_{n}}\left\{\left(\xi_{n}+i 0\right)^{-\frac{\nu}{2}-q-1} \xi_{n}^{-\ell}\left|\xi_{n}\right|^{\nu}\right\} \\
& =\boldsymbol{d}_{1}\left(x^{\prime}\right) \mathscr{F}_{\xi_{n} \rightarrow x_{n}}\left\{\partial_{\xi_{n}}^{m}\left[\left(\xi_{n}+i 0\right)^{\frac{\nu}{2}-q-\ell-1} \theta_{\nu}\left(\xi_{n}\right)\right]\right\}, \tag{B.7.14}
\end{align*}
$$

with $\boldsymbol{d}_{1}\left(x^{\prime}\right):=(-i)^{m} \breve{\mathbf{b}}\left(x^{\prime}\right)\left[\left(i \partial_{x^{\prime}}\right)^{\gamma^{\prime}} \boldsymbol{d}\right]\left(x^{\prime}\right) \in \mathscr{C}^{\infty}\left(\mathbb{R}^{n-1}\right)^{N}$, where we have used the formula, $c f$ Lemma B.5.2,

$$
\begin{equation*}
|t|^{\sigma}=\theta_{\sigma}(t)(t+i 0)^{\sigma} \quad \text { with } \quad \theta_{\sigma}(t)=\chi_{+}(t)+e^{-i \pi \sigma} \chi_{-}(t) \tag{B.7.15}
\end{equation*}
$$

We note that although we have taken derivatives $\partial_{\xi_{n}}^{m}\left|\xi_{n}\right|^{\sigma}$, $\delta$-functions do not appear due to cut off functions (see the beginning of the proof).
Inserting expression (B.7.14) of $\boldsymbol{h}$ into (B.7.12) we find that

$$
\begin{align*}
\boldsymbol{\psi}(x) & =\boldsymbol{p}_{+} \mathscr{F}_{\xi_{n} \rightarrow x_{n}}\left\{\left(\xi_{n}-i\right)^{-\frac{\nu}{2}} \breve{\mathbf{a}}_{0}^{-1}\left(x^{\prime}\right) \mathscr{F}_{x_{n} \rightarrow \xi_{n}}\left[\boldsymbol{h}\left(x^{\prime}, x_{n}\right)\right]\right\} \\
& =\boldsymbol{p}_{+} \mathscr{F}_{\xi_{n} \rightarrow x_{n}}\left\{\left(\xi_{n}-i\right)^{-\frac{\nu}{2}}\left(\xi_{n}+i 0\right)^{\frac{\nu}{2}-q-\ell-m-1} \theta_{\nu}\left(\xi_{n}\right)\right\} \boldsymbol{d}_{2}\left(x^{\prime}\right) \tag{B.7.16}
\end{align*}
$$

with $\boldsymbol{d}_{2}\left(x^{\prime}\right)=c \breve{\mathbf{a}}_{0}^{-1}\left(x^{\prime}\right) \boldsymbol{d}_{1}\left(x^{\prime}\right) \in \mathscr{C}^{\infty}\left(\mathbb{R}^{n-1}\right)^{N}$. By representing the function $\left(\xi_{n}-i\right)^{-\frac{\nu}{2}}$ as a Taylor series in $\left(\xi_{n}-i 0\right)^{-\frac{\nu}{2}-p}, c f$ (B.7.10), and applying the equality

$$
\left(\xi_{n}-i 0\right)^{-\frac{\nu}{2}-p} \theta_{\nu}\left(\xi_{n}\right)=\left(\xi_{n}+i 0\right)^{-\frac{\nu}{2}-p}, \quad p=0,1, \ldots
$$

(see (B.7.15)), we get

$$
\begin{equation*}
\boldsymbol{\psi}(x)=\sum_{p=0}^{M} \mathscr{F}_{\xi_{n} \rightarrow x_{n}}\left\{\left(\xi_{n}+i 0\right)^{-q-\ell-m-p-1}\right\} \boldsymbol{d}_{2, p}\left(x^{\prime}\right)+\boldsymbol{\psi}_{\text {rem }, M+1}(x) \tag{B.7.17}
\end{equation*}
$$

with $\boldsymbol{d}_{2, p} \in \mathscr{C}^{\infty}\left(\mathbb{R}^{n-1}\right)^{N}$ and $\boldsymbol{\psi}_{\text {rem }, M+1} \in \widetilde{\mathrm{H}}^{\left(\infty, s-\frac{\nu}{2}+M+1\right), \infty}\left(\mathbb{R}_{+}^{n}\right)^{N}$. The restriction operator $\boldsymbol{p}_{+}$in front in (B.7.16) was eliminated since the Fourier transform of the analytic function is supported on $\mathbb{R}_{+}$.

From (B.7.9) - (B.7.17) we find

$$
\begin{aligned}
\boldsymbol{\phi}(x) & =\sum_{k=0}^{M-1}\left(D_{n}+i\right)^{-\frac{\nu}{2}} \mathscr{F}_{\xi_{n} \rightarrow x_{n}}\left\{\left(\xi_{n}+i 0\right)^{-k-1}\right\} \boldsymbol{d}_{3, k}\left(x^{\prime}\right)+\left[\boldsymbol{p}_{+} \mathbf{a}^{\infty}\left(x^{\prime} ; D_{n}\right)\right]^{-1} \boldsymbol{g}_{\mathrm{rem}, M+1}^{3} \\
& =\sum_{k=0}^{M-1} \mathscr{F}_{\xi_{n} \rightarrow x_{n}}\left\{\left(\xi_{n}+i\right)^{-\frac{\nu}{2}}\left(\xi_{n}+i 0\right)^{-k-1}\right\} \boldsymbol{d}_{3, k}\left(x^{\prime}\right)+\left[\boldsymbol{p}_{+} \mathbf{a}^{\infty}\left(x^{\prime} ; D_{n}\right)\right]^{-1} \boldsymbol{g}_{\mathrm{rem}, M+1}^{3}
\end{aligned}
$$

By transforming $\left(\xi_{n}+i 0\right)^{-k-1}$ into $\left(\xi_{n}+i\right)^{-k-1}$ as above, and using the asymptotics of $\left[\boldsymbol{p}_{+} \mathrm{a}^{\infty}\left(x^{\prime} ; D_{n}\right)\right]^{-1} \boldsymbol{g}_{\text {rem }, M+1}^{3}$ from Lemma B.6.2, we finally obtain the desired expansion

$$
\begin{equation*}
\boldsymbol{\phi}(x)=\sum_{k=0}^{M-1} x_{n}^{\frac{\nu}{2}+k} e^{-x_{n}} \boldsymbol{d}^{k}\left(x^{\prime}\right)+\boldsymbol{\phi}_{\mathrm{rem}, M+1}(x), \tag{B.7.18}
\end{equation*}
$$

with $\boldsymbol{\phi}_{\mathrm{rem}, M+1} \in \mathrm{H}^{(\infty, s+M+1), \infty}\left(\mathbb{R}_{+}^{n}\right)^{N}$, and $\boldsymbol{d}^{k} \in \mathscr{C}^{\infty}\left(\mathbb{R}^{n-1}\right)^{N}$. The theorem is proved.

## B. 8 Spatial asymptotics of solutions to BVP

We have already described the first two steps of the analysis of asymptotics by the Wiener-Hopf method: (i) the reduction to a $\Psi$ DE (A.3.1) on the boundary, (ii) the asymptotics of the solution of this $\Psi D E$. There remains to derive the spatial asymptotics of the solution $\boldsymbol{u}$ to BVP (A.1.1), represented by the formula (B.2.5), $\boldsymbol{u}=N \boldsymbol{f}+\mathscr{D}[\boldsymbol{u}]-\mathscr{V}[T \boldsymbol{u}]$, if we know the asymptotics of the densities $[\boldsymbol{u}]$ or $[T \boldsymbol{u}]$. Note, that since $\boldsymbol{f} \in \mathscr{C}_{0}^{\infty}\left(\mathbb{R}^{n+1}\right)$, the summand $N \boldsymbol{f}$ only contributes to the regular part of $\boldsymbol{u}$.

Therefore, we only need to apply either the single layer potential $\mathscr{V}$, or the double layer potential $\mathscr{D}$ to a function $\phi$ defined on $\mathscr{M}$, the asymptotic expansion of which being of the form (B.4.3).

Thus, let us by denote by $\mathscr{A}$ either the single layer potential $\mathscr{V}$, or the double layer potential $\mathscr{D}$, see (B.2.6), associated with an homogeneous elliptic second order $N \times N$
system $L\left(D_{\boldsymbol{x}}\right)$ in $\mathbb{R}^{n+1}$ with constant real coefficients ${ }^{(4)}$. Let $q$ be the order of $\mathscr{A}(q=$ -1 if $\mathscr{A}=\mathscr{V}$ and $q=1$ if $\mathscr{A}=\mathscr{D}$ ). We consider $\boldsymbol{u}$ defined on $\Omega$ by

$$
\begin{equation*}
\boldsymbol{u}(\boldsymbol{x})=\mathscr{A} \boldsymbol{\phi}(\boldsymbol{x}), \quad \operatorname{supp} \phi \subset \overline{\mathscr{M}}, \quad \boldsymbol{x} \in \Omega \tag{B.8.1}
\end{equation*}
$$

For any $x^{\prime} \in \mathscr{E}$, let $\tau_{1}\left(\mathscr{X}^{\prime}\right), \ldots, \tau_{\ell}\left(x^{\prime}\right)$ be all different roots of the polynomial equation

$$
\begin{equation*}
\operatorname{det} L\left(\mathscr{I}_{\varkappa}\left(\mathscr{X}^{\prime}\right)(0,1, \tau)\right)=0, \quad \operatorname{Im} \tau<0 \tag{B.8.2}
\end{equation*}
$$

We recall that $(0,1, \tau)$ represents the value of the dual variable $\underline{\xi}=\left(\xi^{\prime}, \xi_{n}, \xi_{n+1}\right)$ and that $\mathscr{I}_{\varkappa}\left(\mathscr{X}^{\prime}\right)$ is the Jacobian of the local coordinate diffeomorphism $\bar{\kappa}, c f$ Definition A.4.1.

We assume that it is possible to enumerate $\tau_{1}\left(\mathscr{X}^{\prime}\right), \ldots, \tau_{\ell}\left(\mathscr{X}^{\prime}\right)$ so that $\left(\mathfrak{H}_{\text {B6 }}\right) \quad$ The multiplicities $n_{1}, \ldots, n_{\ell}$ of $\tau_{1}\left(\mathscr{X}^{\prime}\right), \ldots, \tau_{\ell}\left(\mathscr{X}^{\prime}\right)$ are constant on $\mathscr{E}$.
Therefore the $\tau_{m}$ are $\mathscr{C}^{\infty}(\mathscr{E})$.
Since $L$ is a $N \times N$ elliptic system of order 2 , there holds

$$
n_{1}+\ldots+n_{\ell}=N
$$

and since its coefficients are real, the roots of equation (B.8.2) with $\operatorname{Im} \tau>0$ are the conjugate of the $\tau_{m}\left(\mathscr{X}^{\prime}\right)$. Let for $\mathscr{X}^{\prime} \in \mathscr{E}$ and $m=1, \ldots, \ell$ the angular functions $\psi_{m, \pm}$ be defined as

$$
\begin{equation*}
\psi_{m,-1}\left(\mathscr{X}^{\prime}, \theta\right):=\cos \theta+\tau_{m}\left(\mathscr{X}^{\prime}\right) \sin \theta, \quad \psi_{m,+1}\left(\mathscr{X}^{\prime}, \theta\right):=\cos \theta+\bar{\tau}_{m}\left(\mathscr{X}^{\prime}\right) \sin \theta \tag{B.8.3}
\end{equation*}
$$

Theorem B.8. 1 Let $\phi$ be a $N$-vector function on $\mathscr{M}$ with the following infinite asymptotics without logarithms: $\exists \mu \in \mathbb{R}, \forall K>0$

$$
\boldsymbol{\phi}=\sum_{k=0}^{K-1} r^{\mu+k} \chi(r) \boldsymbol{d}^{k}\left(\mathscr{X}^{\prime}\right)+\boldsymbol{\phi}_{\mathrm{rem}, K}, \quad \boldsymbol{d}^{k} \in \mathscr{C} \mathscr{C}^{\infty}(\mathscr{E})^{N}, \quad \boldsymbol{\phi}_{\mathrm{rem}, K} \in \widetilde{\mathrm{H}}^{\mu+K}(\mathscr{M})^{N}
$$

We assume that $\mu$ is not an integer and that the $N \times N$ second order system $L$ satisfies hypothesis $\left(\mathfrak{H}_{\mathrm{B} 6}\right)$. Let $\mathscr{A}$ denote either the single or the double layer potential associated with $L$ and $q$ its order, and let $\boldsymbol{u}$ be defined on $\Omega$ by $\mathscr{A} \phi$.

Then for arbitrary $K \in \mathbb{N}$, the potential-type function $u$ has in local cylindrical coordinates $\left(\mathscr{X}^{\prime}, r, \theta\right)$ the following asymptotic expansion free of logarithms as well

$$
\begin{align*}
\boldsymbol{u}= & \sum_{m=1}^{\ell} \sum_{\omega= \pm 1} \chi(r)\left[\sum_{j=0}^{n_{m}-1} r^{\mu-q} \sin ^{j} \theta \psi_{m, \omega}^{\mu-q-j}\left(\mathscr{X}^{\prime}, \theta\right) \boldsymbol{d}_{m, \omega}^{j}\left(\mathscr{X}^{\prime}\right)\right.  \tag{B.8.4}\\
& \left.+\sum_{k=1}^{K+q-1} \sum_{j=0}^{p(m, k)} \sum_{|\alpha| \leq N(m, k)} r^{\mu-q+k} \psi_{m, \omega}^{\mu-q-j+k}\left(\mathscr{X}^{\prime}, \theta\right) \sin ^{\alpha_{1}} \theta \cos ^{\alpha_{2}} \theta \boldsymbol{d}_{m, \omega}^{k, j, \alpha}\left(\mathscr{X}^{\prime}\right)\right]+\boldsymbol{u}_{\mathrm{rem}, K}
\end{align*}
$$

where $\boldsymbol{u}_{\mathrm{rem}, K} \in \mathrm{H}_{\mathrm{loc}}^{\mu+K}\left(\mathbb{R}^{n+1}\right)^{N}$ and the coefficients $\boldsymbol{d}_{m, \omega}^{j}$ and $\boldsymbol{d}_{m, \omega}^{j, k, \alpha}$ are $\mathscr{C}^{\infty}(\mathscr{E})$.

[^3]The proof is a direct adaptation of proofs in [ChkDu1, ChkDu2].
As a straightforward corollary of Theorems B.4.2 and B.8.1 combined with

- formulas (B.2.8) and (B.2.10) for Neumann conditions,
- formulas (B.2.9) and (B.2.11) for Dirichlet conditions,
we obtain:
Theorem B.8.2 Let the $N \times N$ second order system L satisfy hypotheses $\left(\mathfrak{H}_{\mathrm{A} 1}\right)$, $\left(\mathfrak{H}_{\mathrm{A} 2}\right)$ and $\left(\mathfrak{H}_{\mathrm{B} 6}\right)$. Then any solution $\boldsymbol{u}$ of BVP (A.1.1) with $\boldsymbol{f} \in \mathscr{C}_{0}^{\infty}\left(\mathbb{R}^{n+1}\right)$ has the following asymptotic expansion in local cylindrical coordinates ( $\mathscr{X}^{\prime}, r, \theta$ )

$$
\begin{align*}
\boldsymbol{u}= & \sum_{m=1}^{\ell} \sum_{\omega= \pm 1} \chi(r)\left[\sum_{j=0}^{n_{m}-1} r^{\frac{1}{2}} \sin ^{j} \theta \psi_{m, \omega}^{\frac{1}{2}-j}\left(x^{\prime}, \theta\right) \boldsymbol{d}_{m, \omega}^{j}\left(\mathscr{X}^{\prime}\right)\right.  \tag{B.8.5}\\
& \left.+\sum_{k=1}^{K-1} \sum_{j=0}^{p(m, k)} \sum_{|\alpha| \leq N(m, k)} r^{\frac{1}{2}+k} \psi_{m, \omega}^{\frac{1}{2}+k-j}\left(x^{\prime}, \theta\right) \sin ^{\alpha_{1}} \theta \cos ^{\alpha_{2}} \theta \boldsymbol{d}_{m, \omega}^{k, j, \alpha}\left(x^{\prime}\right)\right]+\boldsymbol{u}_{\mathrm{rem}, K}
\end{align*}
$$

where $\boldsymbol{u}_{\mathrm{rem}, K} \in \mathrm{H}_{\mathrm{loc}}^{K}\left(\mathbb{R}^{n+1}\right)^{N}$ and the coefficients $\boldsymbol{d}_{m, \omega}^{j}$ and $\boldsymbol{d}_{m, \omega}^{j, k, \alpha}$ are $\mathscr{C}^{\infty}(\mathscr{E})$.
For the particular case of isotropic elasticity we have to deal with the Lamé equation

$$
\begin{equation*}
L\left(D_{\boldsymbol{x}}\right) \boldsymbol{u}=\mu \Delta \boldsymbol{u}+(\lambda+\mu) \operatorname{grad} \operatorname{div} \boldsymbol{u}=\boldsymbol{f}, \quad \boldsymbol{f} \in \mathscr{C}_{0}^{\infty}\left(\mathbb{R}^{3}\right) . \tag{B.8.6}
\end{equation*}
$$

Equation (B.8.2) has one triple root $\tau_{1}=-i$ and for the singular functions (B.8.3) we get

$$
\psi_{1,-1}(\theta)=e^{i \theta} \quad \text { and } \quad \psi_{1,1}(\theta)=e^{-i \theta}
$$

The asymptotics of the displacement $\boldsymbol{u}(\boldsymbol{x})$ has the form

$$
\begin{align*}
& \boldsymbol{u}\left(\mathscr{X}^{\prime}, r, \theta\right)=\sum_{\omega= \pm 1}\left[\sum_{j=0}^{2} r^{\frac{1}{2}} \sin ^{j} \theta e^{i \omega\left(\frac{1}{2}-j\right) \theta} \boldsymbol{d}_{\omega}^{j}\left(\mathscr{X}^{\prime}\right)\right.  \tag{B.8.7}\\
& \left.\quad+\sum_{k=1}^{K-1} \sum_{j=0}^{p_{k}} \sum_{|\alpha| \leq N_{k}} r^{\frac{1}{2}+k} e^{i \omega\left(\frac{1}{2}-j+k\right) \theta} \sin ^{\alpha_{1}} \theta \cos ^{\alpha_{2}} \theta \boldsymbol{d}_{\omega}^{k, j, \alpha}\left(\mathscr{X}^{\prime}\right)\right]+\boldsymbol{u}_{\mathrm{rem}, K}\left(\mathscr{X}^{\prime}, r, \theta\right) .
\end{align*}
$$

The stress $T\left(\boldsymbol{x}, D_{\boldsymbol{x}}\right) \boldsymbol{u}(\boldsymbol{x})$ has a similar asymptotics as the displacement, starting with the exponent $r^{-\frac{1}{2}}$ instead of $r^{\frac{1}{2}}$.

## Part C. The Mellin approach

## C. 1 General edge asymptotics

In our second approach, we are considering the boundary value problem (A.1.1) as a special case of boundary value problems on domains with edges. For such problems, the method of Mellin transformation is a well-developed technique that allows precise descriptions of the solutions in the neighborhood of the edge.

The general description of solutions of problems like (A.1.1) on a wedge originates from Kontratiev's work [Ko1] and was developed in the subsequent works [MaPl1, MaRo1, $\mathrm{NaPl1}]$ and [Da1, CoDa1], among other contributions. As a preparation for our proof on the absence of logarithm, we are going to explain the general edge structure in the framework of the above papers.

We keep the local cylindrical coordinates $\left(\mathscr{X}^{\prime}, r, \theta\right)$ around the edge $\mathscr{E}$, see Definition A.4.1. As this will be of constant use, we introduce the notation $y$ for the two normal cartesian coordinates $\left(x_{n}, \mathscr{X}_{n+1}\right)$, which will be also alternatively denoted by $\left(y_{1}, y_{2}\right)$. Let us consider as domain for the boundary value problem the wedge $W_{\omega}=\mathscr{E} \times \Gamma_{\omega}$ where $\Gamma_{\omega}$ is the plane sector $\{y \sim(r, \theta) \mid \theta \in(-\omega, \omega)\}$ of opening $2 \omega$. Let $\partial_{ \pm} \Gamma_{\omega}$ be the two sides of $\Gamma_{\omega}$. They correspond to the two sides $\partial_{ \pm} W_{\omega}$ of $W_{\omega}$. The situation which is the aim of our investigation corresponds to taking $\omega=\pi$.

But for a while, let us consider the more general case of an elliptic $N \times N$ system $L=\left(L_{k \ell}\right)$ of order $2 d$ complemented by two sets $B_{ \pm}$of $m:=d N$ boundary conditions on $\partial_{ \pm} W_{\omega}$. The general framework of edge asymptotics demands a supplementary condition of ellipticity along the edge, see [MaPl1, MaRo1]. A natural way to satisfy this condition is to suppose that $\left(L, B_{-}, B_{+}\right)$is associated with a coercive form $b$ on $\mathrm{H}^{d}$, see [Da1], as stated in Part A (but now with order $2 d$ and more general boundary conditions).

Thus, let us consider $\boldsymbol{u}$ solution in $\mathrm{H}^{d}\left(W_{\omega}\right)^{N}$ of the following boundary value problem with a right hand side $f \in \mathscr{C}^{\infty}\left(\bar{W}_{\omega}\right)^{N}$

$$
\left\{\begin{align*}
L \boldsymbol{u}=\boldsymbol{f} & \text { in } W_{\omega}  \tag{C.1.1}\\
\gamma_{ \pm} B_{ \pm} \boldsymbol{u}=0 & \text { on } \partial_{ \pm} W_{\omega} .
\end{align*}\right.
$$

The solution $\boldsymbol{u}$ has an infinite edge asymptotics, mainly determined by the expansion of the problem $\left(L, B_{-}, B_{+}\right)$in "homogeneous components" $\left(L^{j}, B_{-}^{j}, B_{+}^{j}\right), j \geq 0$, with respect to the variables $y$ normal to the edge $\mathscr{E}$.

In the coordinates $\left(\mathscr{X}^{\prime}, y\right) \in \mathscr{E} \times \Gamma_{\omega}$, the system $L$ has variable coefficients, in general. We write it with the notation $L=L\left(x^{\prime}, y ; \partial_{x^{\prime}}, \partial_{y}\right)$. For any $x^{\prime} \in \mathscr{E}$, let $L^{0}\left[x^{\prime}\right]$ be the principal part of the operator $L\left(x^{\prime}, 0 ; 0, \partial_{y}\right)$. We denote similarly the boundary operators in local coordinates by $B_{ \pm}\left(\mathscr{X}^{\prime}, r ; \partial_{\mathscr{X}^{\prime}}, \partial_{y}\right)$ and their principal parts in $y=0$ by $B_{ \pm}^{0}\left[\mathscr{X}^{\prime}\right]$.

For each $\mathscr{X}^{\prime}$ fixed in $\mathscr{E}$, the singular exponents associated with $x^{\prime}$ are the complex numbers $\lambda$ such that there exists non-zero solutions $\boldsymbol{\psi}=\boldsymbol{\psi}(\theta)$ to the problem

$$
\left\{\begin{align*}
L^{0}\left[x^{\prime}\right]\left(r^{\lambda} \boldsymbol{\psi}\right)=0 & \text { in } \Gamma_{\omega}  \tag{C.1.2}\\
\gamma_{ \pm} B_{ \pm}^{0}\left[x^{\prime}\right]\left(r^{\lambda} \boldsymbol{\psi}\right)=0 & \text { on } \partial_{ \pm} \Gamma_{\omega} .
\end{align*}\right.
$$

In general, due to the dependency on $\mathscr{X}^{\prime}$ of the coefficients of $\left(L^{0}, B_{ \pm}^{0}\right)$, the set $\Lambda\left[x^{\prime}\right]$ of such $\lambda$ a priori varies with $\mathscr{X}^{\prime} \in \mathscr{E}$, see [MaRo1, CoDa1].

The Ansatz for solutions in the form $r^{\lambda} \boldsymbol{\psi}(\theta)$ has a close relation with the Mellin transform which allows a diagonalization of $\left(L^{0}, B_{ \pm}^{0}\right)\left[x^{\prime}\right]$ for each $x^{\prime}$. Let us recall the Mellin transform $\lambda \mapsto \mathfrak{M}(f)(\lambda)$ of a function $f$ defined on $\mathbb{R}_{+}$:

$$
\mathfrak{M}(f)(\lambda)=\int_{0}^{\infty} r^{\lambda} f(r) \frac{\mathrm{d} r}{r} .
$$

We have the formula $\mathfrak{M}\left(r \partial_{r} f\right)(\lambda)=\lambda \mathfrak{M}(f)(\lambda)$ which is the foundation of the Mellin symbolic calculus. Thus the Mellin symbol $\lambda \mapsto \mathfrak{A}^{0}\left[\mathscr{X}^{\prime}\right](\lambda)$ of problem (C.1.2) is defined after writing $L^{0}$ and $B_{ \pm}^{0}$ in cylindrical coordinates as

$$
r^{-2 d} \mathscr{L}^{0}\left[\mathscr{X}^{\prime}\right]\left(\theta ; r \partial_{r}, \partial_{\theta}\right) \quad \text { and } \quad r^{-\rho_{ \pm, h}} \mathscr{B}_{ \pm, h}^{0}\left[\mathscr{X}^{\prime}\right]\left(\theta ; r \partial_{r}, \partial_{\theta}\right) \quad(h=1, \ldots, m),
$$

(where $\rho_{ \pm, h}=\operatorname{deg} B_{ \pm, h}$ ) by

$$
\begin{aligned}
\mathfrak{A}^{0}\left[\mathscr{X}^{\prime}\right](\lambda): \mathrm{H}^{2 d}(-\omega, \omega)^{N} & \longrightarrow \omega, \omega)^{N} \times \mathbb{C}^{2 m} \\
\varphi & \longmapsto\left(\mathscr{L}^{0}\left[\mathscr{x}^{\prime}\right]\left(\lambda, \partial_{\theta}\right) \varphi, \gamma_{ \pm} \mathscr{B}_{ \pm}^{0}\left[\mathscr{X}^{\prime}\right]\left(\lambda, \partial_{\theta}\right) \varphi\right) .
\end{aligned}
$$

For each $\mathscr{X}^{\prime} \in \mathscr{E}, \lambda \mapsto \mathfrak{A}^{0}\left[\mathscr{X}^{\prime}\right](\lambda)^{-1}$ is meromorphic in $\mathbb{C}$ and the set of its poles is $\Lambda\left[\mathscr{X}^{\prime}\right]$.
It is possible to classify the singularities occurring in the asymptotics of a solution $\boldsymbol{u}$ of (C.1.1) in (i) Leading singularities and (ii) Shadow singularities.
(i) The leading singularities $s^{0}$ of $\boldsymbol{u}$ are directly obtained from the Mellin transform $\lambda \mapsto \mathfrak{M}(\boldsymbol{f})\left[\mathscr{X}^{\prime}\right](\lambda)$ of $\boldsymbol{f}^{(5)}$ via the Mellin symbol $\mathfrak{A}^{0}\left[\mathscr{X}^{\prime}\right]$ of problem (C.1.2) by the inverse Mellin formula

$$
\begin{equation*}
s^{0}\left(\mathscr{X}^{\prime}, y\right)=\frac{1}{2 i \pi} \int_{\gamma^{0}} r^{\mu}\left[\mathfrak{A}^{0}\left[\mathscr{X}^{\prime}\right](\mu)\right]^{-1} \mathfrak{M}\left(r^{2 d} \boldsymbol{f}, 0,0\right)\left[\mathscr{X}^{\prime}\right](\mu) \mathrm{d} \mu, \tag{C.1.3}
\end{equation*}
$$

where the 0 in $(\boldsymbol{f}, 0,0)$ stand for the zero boundary conditions and $\gamma^{0}$ is a suitable contour surrounding the poles $\lambda \in \Lambda\left[\mathscr{X}^{\prime}\right]$ in the right half plane $\operatorname{Re} \lambda>d-1 .^{(6)}$

[^4](ii) The shadow singularities require for their definition the Taylor expansion of the coefficients of $L$ and $B_{ \pm}$with respect to $y$ : Let
$$
L=\sum_{|i|+|k| \leq 2 d} \ell^{i, k}\left(\mathscr{X}^{\prime}, y\right) \partial_{\mathscr{X}}^{i} \partial_{y}^{k} \quad \text { and } \quad B_{ \pm, h}=\sum_{|i|+|k| \leq \rho_{ \pm, h}} b_{ \pm, h}^{i, k}\left(\mathscr{X}^{\prime}, y\right) \partial_{\mathscr{X}}^{i} \partial_{y}^{k}
$$
be the expressions of $L$ and $B_{ \pm}$. Then for $j \in \mathbb{N}$ we define
\[

$$
\begin{aligned}
L^{j}\left[\mathscr{X}^{\prime}\right] & :=\sum_{|i| \leq 2 d} \sum_{|k|-|\beta|=2 d-j} \partial_{y}^{\beta} \ell^{i, k}\left(\mathscr{X}^{\prime}, 0\right) \frac{y^{\beta}}{\beta!} \partial_{\mathscr{X}^{\prime}}^{i} \partial_{y}^{k} \\
B_{ \pm, h}^{j}\left[\mathscr{X}^{\prime}\right] & :=\sum_{|i| \leq \rho_{ \pm, h}} \sum_{|k|-|\beta|=\rho_{ \pm, h}-j} \partial_{y}^{\beta} b_{ \pm, h}^{i, k}\left(\mathscr{X}^{\prime}, 0\right) \frac{y^{\beta}}{\beta!} \partial_{\mathscr{X}^{\prime}}^{i} \partial_{y}^{k} .
\end{aligned}
$$
\]

Let $\mathfrak{A}^{j}\left[\mathscr{X}^{\prime}\right]$ denote the triple $\left(L^{j}\left[\mathscr{X}^{\prime}\right], B_{ \pm}^{j}\left[\mathscr{X}^{\prime}\right]\right)$. Then the shadow singularities $s^{1}, \ldots$, $s^{p}$ are recursively defined as

$$
\begin{gathered}
\boldsymbol{s}^{p}\left(\mathscr{X}^{\prime}, y\right)=-\frac{1}{2 i \pi} \int_{\gamma^{0}+p} r^{\mu}\left[\mathfrak{A}^{0}\left[\mathscr{X}^{\prime}\right](\mu)\right]^{-1} \mathfrak{M}\left(r^{\boldsymbol{\beta}}\left(\boldsymbol{f}^{p}, \boldsymbol{g}_{ \pm}^{p}\right)\right)\left[\mathscr{X}^{\prime}\right](\mu) \mathrm{d} \mu, \\
\text { with } \quad\left(\boldsymbol{f}^{p}, \boldsymbol{g}_{ \pm}^{p}\right)=\mathfrak{A}^{1} \boldsymbol{s}^{p-1}+\cdots+\mathfrak{A}^{p} \boldsymbol{s}^{0} .
\end{gathered}
$$

Here $\boldsymbol{\beta}$ is the collection of degrees $\left(2 d, \ldots, 2 d, \rho_{ \pm, 1}, \ldots, \rho_{ \pm, m}\right)$ and $r^{\boldsymbol{\beta}}\left(\boldsymbol{f}, \boldsymbol{g}_{ \pm}\right)$is a condensed notation for ( $r^{2 d} \boldsymbol{f}, r^{\rho_{-, 1}} g_{-, 1}, \ldots, r^{\rho_{-, m}} g_{-, m}, r^{\rho_{+, 1}} g_{+, 1}, \ldots, r^{\rho_{+, m}} g_{+, m}$ ).

Then the sum $s^{0}+s^{1}+\ldots+s^{p}+\ldots$ gives the asymptotics of $\boldsymbol{u}$ as $r \rightarrow 0$.
In the most general case, the structure of the $s^{p}$ is quite difficult to describe because of the possible change of multiplicities in the singular exponents $\lambda\left[x^{\prime}\right]$, see [CoDa1, CoDa3]. If hypotheses are made to avoid any change of multiplicity, see [MaRo1], each $s^{p}$ can be decomposed into elementary terms of the form $c\left(x^{\prime}\right) r^{\lambda\left(x^{\prime}\right)+p} \log ^{q} r \varphi\left(x^{\prime}, \theta\right)$. Thus we obtain the following expansion in local cylindrical coordinates: For any $K \in \mathbb{N}$

$$
\begin{equation*}
\boldsymbol{u}=\sum_{\operatorname{Re} \lambda+p \leq K} \sum_{q=0}^{q(\lambda, p)} \sum_{j=1}^{j(\lambda, p, q)} c_{j}^{\lambda, p, q}\left(\mathscr{X}^{\prime}\right) r^{\lambda\left(\mathscr{X}^{\prime}\right)+p} \log ^{q} r \boldsymbol{\varphi}_{j}^{\lambda, p, q}\left(\mathscr{X}^{\prime}, \theta\right)+\boldsymbol{u}_{\mathrm{rem}, K} \tag{C.1.5}
\end{equation*}
$$

The exponents $\lambda\left(x^{\prime}\right)$ belong to $\Lambda\left[x^{\prime}\right] \cup \mathbb{N}$ and their real part is $>d-1$. The coefficients $c_{j}^{\lambda, p, q}$ are $\mathscr{C}^{\infty}$ functions on $\mathscr{E}$ and depend on $\boldsymbol{f}$. The remainder $\boldsymbol{u}_{\text {rem }, K}$ satisfies $\partial^{\beta} \boldsymbol{u}_{\mathrm{rem}, K}=o\left(r^{K-|\beta|+1 / 2}\right)$ as $r \rightarrow 0$ for any multi-index $\beta \in \mathbb{N}_{0}^{n+1}$. The $\boldsymbol{\varphi}_{j}^{\lambda, p, q}$ are angular $N$-component vector functions in $\mathscr{C}^{\infty}([-\omega, \omega] \times \mathscr{E})$ and depend only on the domain $\Omega$ and the operators $(L, B)$.

The $\log r$ terms come either from non-trivial Jordan chains in $\mathfrak{A}^{0}\left[\mathscr{x}^{\prime}\right]^{-1}$, or from resonances between $\mathfrak{A}^{0}\left[\mathscr{X}^{\prime}\right]^{-1}$ and the Mellin transforms $\mathfrak{M}\left(r^{2 d} \boldsymbol{f}, 0,0\right)\left[\mathscr{X}^{\prime}\right]$, see (C.1.3), or $\mathfrak{M}\left(r^{\boldsymbol{\beta}} \mathfrak{A}^{1} s^{p-1}+\cdots\right)\left[\mathscr{X}^{\prime}\right]$, see (C.1.4).

## C. 2 Crack asymptotics, First results

From now on, we concentrate on the situation of a crack, i.e. when the opening $\omega$ is $\pi$, and when the same boundary conditions are applied on both sides of the crack, i.e. $B_{ \pm}=B$. Thus the boundary conditions are denoted by $B=\left(B_{1}, \ldots, B_{m}\right)$ and the order of $B_{h}$ is $\rho_{h}, h=1, \ldots, m$. The boundary problem takes then the form

$$
\left\{\begin{align*}
L \boldsymbol{u}=\boldsymbol{f} & \text { in } W_{\pi}  \tag{C.2.1}\\
\gamma_{ \pm} B \boldsymbol{u}=0 & \text { on } \partial_{ \pm} W_{\pi}
\end{align*}\right.
$$

where we assume that $f \in \mathscr{C}_{0}^{\infty}\left(\mathbb{R}^{n+1}\right)$.
In this situation there holds

$$
\begin{equation*}
\forall \mathscr{X}^{\prime} \in \mathscr{E}, \quad \Lambda\left[x^{\prime}\right]=\left\{\frac{k}{2} ; k \in \mathbb{Z}\right\} . \tag{C.2.2}
\end{equation*}
$$

This has been known for a long time for the Laplace operator, see [Gr1]. It is proved for elasticity systems in [DuWe1], for general second order Petrovskii-elliptic systems (such as thermoelasticity or electroelasticity for example) in [ChkDu1, ChkDu2], for general scalar elliptic Dirichlet problems of order $2 m$ in [Koz1], and finally in the general framework of Agmon-Douglis-Nirenberg elliptic systems in [CoDa4].

Therefore the assumptions on the constant multiplicity of the singular exponents are satisfied and expansion (C.1.5) holds with $\lambda\left(x^{\prime}\right)=\frac{k}{2}$. This clear separation of the spectrum allows a decomposition of leading singularity $s^{0}$ in (quasi-)homogeneous elementary parts $\Phi_{\lambda}^{0}$ for $\lambda$ of the form $\lambda=\frac{k}{2}$ according to:

$$
\begin{equation*}
\Phi_{\lambda}^{0}\left[\mathscr{x}^{\prime}\right]=\frac{1}{2 i \pi} \int_{\gamma(\lambda)} r^{\mu} \mathfrak{A}^{0}\left[\mathscr{X}^{\prime}\right](\mu)^{-1} \mathfrak{M}\left(r^{2 d} \boldsymbol{f}, 0,0\right)\left[\mathscr{X}^{\prime}\right](\mu) \mathrm{d} \mu \tag{C.2.3}
\end{equation*}
$$

where $\gamma(\lambda)$ is the circle with center $\lambda$ and radius $\frac{1}{4}$.
Definition C.2.1 If $\Phi^{0}$ is defined by a residue formula like (C.2.3) on the circle $\gamma(\lambda)$, we call sequence of shadows associated with $\Phi^{0}$, the infinite sequence $\Phi^{p}, p \geq 1$, defined by

$$
\begin{equation*}
\Phi^{p}\left[\mathscr{X}^{\prime}\right]=-\frac{1}{2 i \pi} \int_{\gamma(\lambda)+p} r^{\mu} \mathfrak{A}^{0}\left[\mathscr{X}^{\prime}\right](\mu)^{-1} \mathfrak{M}\left(r^{\boldsymbol{\beta}} \mathbb{1}_{r \in[0,1]}\left(\mathfrak{A}^{1} \Phi^{p-1}+\cdots+\mathfrak{A}^{p} \Phi^{0}\right)\right)\left[\mathscr{X}^{\prime}\right](\mu) \mathrm{d} \mu \tag{C.2.4}
\end{equation*}
$$

Here $\gamma(\lambda)+p$ is the contour around $\lambda+p$ translated from $\gamma(\lambda)$ and $\mathbb{1}_{r \in[0,1]}$ is the characteristic function in $r$ of the interval $[0,1]$.

By linearity, we obtain that a decomposition of $s^{0}$ in a sum of $\Phi_{\lambda}^{0}$ provides the corresponding decomposition of the shadow $s^{p}$ in a sum $\Phi_{\lambda}^{p}$, where $\left(\Phi^{p}\right)_{p}$ is the sequence of shadows associated with $\Phi_{\lambda}^{0}$. Therefore, from now on we only consider elementary leading singularities of the form (C.2.3) and their sequence of shadows.

The result of [CoDa4, Thms $5.2 \& 5.3]$ gives moreover:
(i) In the leading singularities the non-integer exponents $k+\frac{1}{2}$ have no $\log r$ terms and the corresponding basis of singular functions $\left(r^{k+1 / 2} \boldsymbol{\varphi}_{j}^{k+1 / 2}\right)_{j}$ has the dimension $m$.
(ii) Let $\rho_{\max }:=\max \left\{\rho_{1}, \ldots, \rho_{m}\right\}$. For any integer $\lambda \geq \rho_{\max }$, the functions $r^{\lambda} \boldsymbol{\varphi}_{j}^{\lambda}(\theta)$ are polynomials in the variables $\left(y_{1}, y_{2}\right)$. Moreover the shadows of polynomials are polynomials.
Therefore:
(i) For an exponent $\lambda=\frac{1}{2}+k$, the elementary leading singularities have the form

$$
\begin{equation*}
\Phi_{\lambda}^{0}\left[\mathscr{X}^{\prime}\right]=\sum_{j=1}^{m} c_{j}\left(\mathscr{X}^{\prime}\right) r^{k+1 / 2} \boldsymbol{\varphi}_{j}^{k+1 / 2}\left(\mathscr{X}^{\prime}, \theta\right) \quad c_{j} \in \mathscr{C}^{\infty}(\mathscr{E}) . \tag{C.2.5}
\end{equation*}
$$

(ii) For a positive integer $\lambda \geq \rho_{\max }, \Phi_{\lambda}^{0}$ is a finite sum of terms of the form $c\left(x^{\prime}\right) \psi(y)$ with smooth $c$ and polynomial $\psi$ (homogeneous of degree $\lambda$ ). Moreover, the sequence of shadows $\Phi_{\lambda}^{p}$ associated with $\Phi_{\lambda}^{0}$ have a similar structure with homogeneous polynomials of degree $\lambda+p$.
As a consequence, we have obtained the statement of Proposition A.4.2.
But, when $\lambda=\frac{1}{2}+k$, since $\lambda+p=\frac{1}{2}+k+p$ is a singular exponent, i.e. a pole of $\mathfrak{A}^{0}\left[\mathscr{X}^{\prime}\right]^{-1}$, we should expect resonances inside the integrand of the shadows $\Phi_{\lambda}^{p}$, between $\mathfrak{M}\left(r^{\boldsymbol{\beta}} \mathfrak{A}^{1} \Phi^{p-1}+\cdots\right)\left[\mathscr{x}^{\prime}\right]$ and $\mathfrak{A}^{0}\left[\mathscr{X}^{\prime}\right]^{-1}$, i.e. poles of order $>1$, which would yield $\log r$ factors. We are going to prove that, in fact, there are no resonances.

## C. 3 "CAYLEY" REPRESENTATION FORMULAE

Our method is a direct continuation of [CoDa2] where "Cayley representation formulae" are introduced to describe the angular behavior (in $\theta$ ) of the singular functions. Its is shown there that any singularity can be expressed by combination of two fundamental types of functions which, using the complex writing $\zeta$ of the cartesian variables $y=\left(y_{1}, y_{2}\right)$

$$
\zeta=y_{1}+i y_{2}=r e^{i \theta}
$$

can be written as, for any $\lambda \in \mathbb{C}, \zeta \in \mathbb{C}$ with $\zeta \notin \mathbb{R}^{-}$, and $\alpha \in \mathbb{C}$ with $|\alpha|<1$ :

$$
(\alpha \zeta+\bar{\zeta})^{\lambda} \quad \text { and } \quad(\zeta+\alpha \bar{\zeta})^{\lambda}
$$

The above functions have to be interpreted in the following way:

$$
\begin{equation*}
(\alpha \zeta+\bar{\zeta})^{\lambda}=\bar{\zeta}^{\lambda}\left(1+\alpha \frac{\zeta}{\bar{\zeta}}\right)^{\lambda} \quad \text { and } \quad(\zeta+\alpha \bar{\zeta})^{\lambda}=\zeta^{\lambda}\left(1+\alpha \frac{\bar{\zeta}}{\zeta}\right)^{\lambda} \tag{C.3.1}
\end{equation*}
$$

which means in polar coordinates $r>0, \theta \in(-\pi, \pi)$ :

$$
\begin{equation*}
(\alpha \zeta+\bar{\zeta})^{\lambda}=r^{\lambda} e^{-i \theta \lambda}\left(1+\alpha e^{2 i \theta}\right)^{\lambda} \quad \text { and } \quad(\zeta+\alpha \bar{\zeta})^{\lambda}=r^{\lambda} e^{i \theta \lambda}\left(1+\alpha e^{-2 i \theta}\right)^{\lambda} \tag{C.3.2}
\end{equation*}
$$

The action of a partial differential operator $Q\left(\partial_{1}, \partial_{2}\right)$ on $(\alpha \zeta+\bar{\zeta})^{\lambda}$ and $(\zeta+\alpha \bar{\zeta})^{\lambda}$ exhibits its Cayley symbols $Q^{+}(\alpha)$ and $Q^{-}(\alpha)$ as follows:

$$
Q^{+}(\alpha):=Q(\alpha+1, i(\alpha-1)) \quad \text { and } \quad Q^{-}(\alpha):=Q(1+\alpha, i(1-\alpha))
$$

and there holds, if $Q$ is homogeneous of degree $q$

$$
\left\{\begin{array}{l}
Q\left(\partial_{y}\right)(\alpha \zeta+\bar{\zeta})^{\lambda}=P_{q}(\lambda)(\alpha \zeta+\bar{\zeta})^{\lambda-q} Q^{+}(\alpha) \\
Q\left(\partial_{y}\right)(\zeta+\alpha \bar{\zeta})^{\lambda}=P_{q}(\lambda)(\zeta+\alpha \bar{\zeta})^{\lambda-q} Q^{-}(\alpha)
\end{array}\right.
$$

where $P_{q}(\lambda)$ is the polynomial of degree $q, \lambda(\lambda-1) \cdots(\lambda-q+1)$.
Let us fix $x^{\prime} \in \mathscr{E}$. Let $L^{ \pm}\left[x^{\prime}\right](\alpha)$ be the two Cayley symbols of $L^{0}\left[x^{\prime}\right]$ and $B^{ \pm}\left[x^{\prime}\right]$ those of $B^{0}\left[\mathscr{X}^{\prime}\right]$. We have the following formulas, valid for any $\alpha \in \mathbb{C}$, which are the matrix version of the above ones: let $\boldsymbol{q} \in \mathbb{C}^{N}$ be a vector, there holds

$$
\left\{\begin{align*}
L^{0}\left[\mathscr{X}^{\prime}\right]\left(\partial_{y}\right)\left\{(\alpha \zeta+\bar{\zeta})^{\lambda} \boldsymbol{q}\right\} & =P_{2 d}(\lambda)(\alpha \zeta+\bar{\zeta})^{\lambda-2 d} L^{+}\left[\mathscr{X}^{\prime}\right](\alpha) \boldsymbol{q}  \tag{C.3.3}\\
L^{0}\left[\mathscr{x}^{\prime}\right]\left(\partial_{y}\right)\left\{(\zeta+\alpha \bar{\zeta})^{\lambda} \boldsymbol{q}\right\} & =P_{2 d}(\lambda)(\zeta+\alpha \bar{\zeta})^{\lambda-2 d} L^{-}\left[\mathscr{x}^{\prime}\right](\alpha) \boldsymbol{q}
\end{align*}\right.
$$

These Cayley symbols allow to describe for any $\mathscr{X}^{\prime}$ and $\lambda$ the space $\mathfrak{Z}\left[\mathscr{X}^{\prime}\right](\lambda)$ of the homogeneous functions $\boldsymbol{v}$ of degree $\lambda$, solutions of the equation without boundary conditions

$$
L^{0}\left[x^{\prime}\right] \boldsymbol{v}=0 .
$$

Due to the ellipticity of the operator $L^{0}\left[x^{\prime}\right]$, the equations $\operatorname{det} L^{ \pm}\left[x^{\prime}\right](\alpha)=0$ have $m$ roots inside the unit disc $|\alpha|<1$, counting multiplicity, and no roots on the unit circle $|\alpha|=1$. Let us denote

$$
\begin{equation*}
\alpha_{1}^{-}\left[x^{\prime}\right], \ldots, \alpha_{m_{-}}^{-}\left[x^{\prime}\right], \quad \alpha_{1}^{+}\left[x^{\prime}\right], \ldots, \alpha_{m_{+}}^{+}\left[\mathscr{x}^{\prime}\right] \tag{C.3.4}
\end{equation*}
$$

the distinct roots of $\operatorname{det} L^{-}\left[x^{\prime}\right]$ and $\operatorname{det} L^{+}\left[x^{\prime}\right]$ inside the unit disc.
For a while let us assume that these roots are simple (i.e. $m_{ \pm}=m$ ). Thus, let $\boldsymbol{q}_{\ell}^{ \pm}\left[\mathscr{X}^{\prime}\right] \in$ $\mathbb{C}^{N}$ be non-zero elements of the kernels $\operatorname{ker} L^{ \pm}\left(\alpha_{\ell}^{ \pm}\right)$, and for any (non-integer) $\lambda \in \mathbb{C}$ let us define the $N$-component functions

$$
\boldsymbol{w}_{\ell}^{+}\left[\mathscr{x}^{\prime}\right](\lambda):=\left(\alpha_{\ell}^{+}\left[\mathscr{x}^{\prime}\right] \zeta+\bar{\zeta}\right)^{\lambda} \boldsymbol{q}_{\ell}^{+}\left[\mathscr{x}^{\prime}\right] \quad \text { and } \quad \boldsymbol{w}_{\ell}^{-}\left[\mathscr{x}^{\prime}\right](\lambda):=\left(\zeta+\alpha_{\ell}^{-}\left[x^{\prime}\right] \bar{\zeta}\right)^{\lambda} \boldsymbol{q}_{\ell}^{-}\left[x^{\prime}\right] .
$$

Formulas (C.3.3) give immediately that these functions solve the equation $L^{0}\left[x^{\prime}\right] \boldsymbol{v}=0$, thus belong to $\mathfrak{Z}\left[x^{\prime}\right](\lambda)$.

As proved in [CoDa2, Th. 2.1], these $2 m$ functions form a basis of the space $\mathfrak{Z}\left[\mathscr{x}^{\prime}\right](\lambda)$ and, moreover, we obtain "stable" expressions of $\boldsymbol{w}_{\ell}^{ \pm}\left[\mathscr{x}^{\prime}\right](\lambda)$ with respect to the parameter $x^{\prime}$ without the assumptions that the roots $\alpha_{\ell}^{+}\left[x^{\prime}\right]$ are simple, by using contour integrals in $\alpha$ around the disc $D_{\delta}$ of radius with $\delta<1$ such that $D_{\delta}$ contains all roots $\alpha_{\ell}^{ \pm}\left[\mathscr{X}^{\prime}\right]$ :

There exists $N$-component polynomials of degree $d-1$ in $\alpha$ depending smoothly on $\mathscr{X}^{\prime}$, denoted $\boldsymbol{q}_{\ell}^{ \pm}\left[\mathscr{x}^{\prime}\right](\alpha)$ for $\ell=1, \ldots, m$, which define a basis $\left\{\boldsymbol{w}_{\ell}^{ \pm}\left[\mathscr{x}^{\prime}\right]\right\}$ of $\mathfrak{Z}\left[x^{\prime}\right](\lambda)$ :

$$
\left\{\begin{align*}
\boldsymbol{w}_{\ell}^{+}\left[\mathscr{x}^{\prime}\right](\lambda) & =\int_{|\alpha|=\delta}(\alpha \zeta+\bar{\zeta})^{\lambda} L^{+}\left[\mathscr{x}^{\prime}\right](\alpha)^{-1} \boldsymbol{q}_{\ell}^{+}\left[\mathscr{x}^{\prime}\right](\alpha) \mathrm{d} \alpha  \tag{C.3.5}\\
\boldsymbol{w}_{\ell}^{-}\left[\mathscr{X}^{\prime}\right](\lambda) & =\int_{|\alpha|=\delta}(\zeta+\alpha \bar{\zeta})^{\lambda} L^{-}\left[\mathscr{x}^{\prime}\right](\alpha)^{-1} \boldsymbol{q}_{\ell}^{-}\left[\mathscr{x}^{\prime}\right](\alpha) \mathrm{d} \alpha
\end{align*}\right.
$$

This basis allows the construction of a $2 m \times 2 m$ matrix $\mathscr{N}\left[\mathscr{x}^{\prime}\right](\lambda)$ whose inverse has the same poles as the inverse of the Mellin symbol $\mathfrak{A}\left[\mathscr{x}^{\prime}\right](\lambda)^{-1}$ : For this let us introduce $\mathfrak{W}\left[x^{\prime}\right](\lambda)$ the $N \times 2 m$ matrix the $2 m$ columns of which are

$$
\boldsymbol{w}_{1}^{+}\left[\mathscr{x}^{\prime}\right](\lambda), \ldots, \boldsymbol{w}_{m}^{+}\left[\mathscr{x}^{\prime}\right](\lambda), \boldsymbol{w}_{1}^{-}\left[x^{\prime}\right](\lambda), \ldots, \boldsymbol{w}_{m}^{-}\left[\mathscr{x}^{\prime}\right](\lambda) .
$$

Let us recall that $B^{0}\left[x^{\prime}\right]$ is the $m \times N$ matrix of the principal parts of the boundary operators $B\left(\mathscr{X}^{\prime}, 0 ; 0, \partial_{y}\right)$. Let $\mathfrak{g}_{ \pm}$be the trace operators (acting on homogeneous functions) ${ }^{(7)}$

$$
\mathfrak{g}_{-} v=\left.v\right|_{r=1 \text { and } \theta=-\pi} \quad \text { and } \quad \mathfrak{g}_{+} v=\left.v\right|_{r=1 \text { and } \theta=\pi} .
$$

The characteristic matrix of the problem is then the $2 m \times 2 m$ scalar matrix given by

$$
\mathscr{N}\left[\mathscr{X}^{\prime}\right](\lambda)=\binom{\mathfrak{g}_{-} B^{0}\left[\mathscr{X}^{\prime}\right]}{\mathfrak{g}_{+} B^{0}\left[\mathscr{x}^{\prime}\right]} \mathfrak{W}\left[\mathscr{x}^{\prime}\right](\lambda) .
$$

The formula describing $\left[\mathfrak{A}^{0}\left[\mathscr{x}^{\prime}\right](\mu)\right]^{-1}$ involves a right inverse to the operator $L^{0}$ on homogeneous functions of degree $\lambda$ (i.e. without boundary conditions) and the inverse of the matrix $\mathscr{N}\left[\mathscr{X}^{\prime}\right](\lambda)$ allows the correction of boundary conditions.

Let $\mathfrak{H}^{\lambda}$ be the space of $N$-component vector functions homogeneous of degree $\lambda$ on the plane sector $\Gamma_{\pi}$. And let $\boldsymbol{f} \mapsto \boldsymbol{v}=\mathfrak{R}\left[\mathscr{x}^{\prime}\right](\lambda) \boldsymbol{f}$ be a solution operator of the problem $L^{0}\left[\mathscr{X}^{\prime}\right] \boldsymbol{v}=\boldsymbol{f}$, acting from $\mathfrak{H}^{\lambda-2 d}$ into $\mathfrak{H}^{\lambda}$. According to [CoDa2], it is possible to construct such an operator with $\mathscr{C}^{\infty}$ regularity in $\mathscr{X}^{\prime}$ and analytic dependency in $\lambda$.

Our first representation theorem for the inverse symbol $\mathfrak{A}^{0}\left[\mathscr{X}^{\prime}\right]^{-1}$ is the following, see [CoDa2, Th. 4.4], - We write it directly for the Mellin integrand $r^{\mu}\left[\mathfrak{A}^{0}\left[\mathscr{X}^{\prime}\right](\mu)\right]^{-1}$ in view of application in formulas (C.2.3) and (C.2.4):
Theorem C.3.1 Let $\mathfrak{R}\left[\mathscr{X}^{\prime}\right](\lambda)$ be a right inverse to $L^{0}\left[\mathscr{x}^{\prime}\right]$, acting from $\mathfrak{H}^{\lambda-2 d}$ into $\mathfrak{H}^{\lambda}$. We have for any $\mathscr{X}^{\prime} \in \mathscr{E}, \mu \in \mathbb{C}$ and any $\left(\boldsymbol{F}, \boldsymbol{G}_{ \pm}\right) \in \mathrm{L}^{2}(-\pi, \pi)^{N} \times \mathbb{C}^{m} \times \mathbb{C}^{m}$ :

$$
\begin{align*}
r^{\mu}\left[\mathfrak{A}^{0}\left[\mathscr{X}^{\prime}\right](\mu)\right]^{-1} & \left(\boldsymbol{F}, \boldsymbol{G}_{ \pm}\right)= \\
& \mathfrak{R}\left[\mathscr{X}^{\prime}\right](\mu)\left(r^{\mu-2 d} \boldsymbol{F}\right)+  \tag{C.3.6}\\
& \mathfrak{W}\left[\mathscr{X}^{\prime}\right](\mu) \mathscr{N}\left[\mathscr{X}^{\prime}\right](\mu)^{-1}\left(\boldsymbol{G}_{ \pm}-\mathfrak{g}_{ \pm} B^{0}\left[\mathscr{X}^{\prime}\right] \mathfrak{R}\left[\mathscr{X}^{\prime}\right](\mu)\left(r^{\mu-2 d} \boldsymbol{F}\right)\right) .
\end{align*}
$$

[^5]Formula (C.3.6) will be applied recursively to special subsets of triples ( $\boldsymbol{F}, \boldsymbol{G}_{ \pm}$) which have the property to be the traces (in $r=1$ ) of homogeneous functions representable by Cayley integrals like (C.3.5):

Definition C.3.2 For any $\lambda \in \mathbb{C}$, let us denote by $\mathfrak{H}_{0}^{\lambda}$ the subspace of homogeneous $N$ component functions $f \in \mathfrak{H}^{\lambda}$ which admit a representation as:

$$
\begin{equation*}
\boldsymbol{f}=\int_{|\alpha|=\delta}(\alpha \zeta+\bar{\zeta})^{\lambda} \boldsymbol{q}^{+}(\alpha) \mathrm{d} \alpha+\int_{|\alpha|=\delta}(\zeta+\alpha \bar{\zeta})^{\lambda} \boldsymbol{q}^{-}(\alpha) \mathrm{d} \alpha \tag{C.3.7}
\end{equation*}
$$

with $N$-component vectors $\boldsymbol{q}^{ \pm}$meromorphic in $\alpha$ (and without pole in the annulus $\delta \leq$ $|\alpha| \leq 1$ ). Such a representation is made unique if we assume that the $\boldsymbol{q}^{ \pm}$are holomorphic outside the unit disc and tend to 0 as $|\alpha| \rightarrow \infty$.

We can define a special solution operator $\mathfrak{R}_{0}\left[\mathscr{X}^{\prime}\right](\lambda)$ acting on the subspace $\mathfrak{H}_{0}^{\lambda-2 d}$ into $\mathfrak{H}_{0}^{\lambda}:$ For $\boldsymbol{f} \in \mathfrak{H}_{0}^{\lambda-2 d}$ represented by (C.3.7) with the uniqueness constraint, we define $\mathfrak{R}_{0}\left[\mathscr{X}^{\prime}\right](\lambda) \boldsymbol{f}$ by

$$
\begin{align*}
& \mathfrak{R}_{0}\left[\mathscr{X}^{\prime}\right](\lambda) \boldsymbol{f}=P_{2 d}(\lambda)^{-1} \int_{|\alpha|=\delta}(\alpha \zeta+\bar{\zeta})^{\lambda} L^{+}\left[\mathscr{X}^{\prime}\right](\alpha)^{-1} \boldsymbol{q}^{+}(\alpha) \mathrm{d} \alpha \\
&+P_{2 d}(\lambda)^{-1} \int_{|\alpha|=\delta}(\zeta+\alpha \bar{\zeta})^{\lambda} L^{-}\left[\mathscr{X}^{\prime}\right](\alpha)^{-1} \boldsymbol{q}^{-}(\alpha) \mathrm{d} \alpha . \tag{C.3.8}
\end{align*}
$$

The vector function obviously belongs to $\mathfrak{H}_{0}^{\lambda}$ and if $P_{2 d}(\lambda) \neq 0$, formulae (C.3.3) give immediately that $L^{0}\left[\mathscr{X}^{\prime}\right] \Re_{0}\left[\mathscr{X}^{\prime}\right](\lambda) \boldsymbol{f}=\boldsymbol{f}$.

## C. 4 Representation of singularities

We start from the expression (C.2.3) of the leading singularity $\Phi^{0}$. The function

$$
\left(\mathscr{X}^{\prime}, \mu\right) \quad \mapsto \quad \mathfrak{M}\left(r^{2 d} \boldsymbol{f}, 0,0\right)\left[x^{\prime}\right](\mu)
$$

is $\mathscr{C}^{\infty}(\mathscr{E})$ in $\mathscr{X}^{\prime}$ and analytic in $\mu$ in the disc $\delta_{\lambda}$ encircled by the contour $\gamma(\lambda)$. Using the representation (C.3.6) with analytic $F$ and zero $G_{ \pm}$, we find that the only pole of $r^{\mu}\left[\mathfrak{A}^{0}\left[\mathscr{X}^{\prime}\right](\mu)\right]^{-1} \mathfrak{M}\left(r^{2 d} \boldsymbol{f}, 0,0\right)\left[\mathscr{X}^{\prime}\right](\mu)$ inside $\delta_{\lambda}$ is $\mu=\lambda$ and that there holds

$$
\begin{equation*}
\Phi^{0}\left[\mathscr{x}^{\prime}\right]=\frac{1}{2 i \pi} \int_{\gamma(\lambda)} \mathfrak{W}\left[\mathscr{x}^{\prime}\right](\mu) \mathscr{N}\left[\mathscr{x}^{\prime}\right](\mu)^{-1} \boldsymbol{\psi}^{0}\left[\mathscr{x}^{\prime}\right](\mu) \mathrm{d} \mu \tag{C.4.1}
\end{equation*}
$$

with a $2 m$-component vector function $\left(x^{\prime}, \mu\right) \mapsto \boldsymbol{\psi}^{0}\left[x^{\prime}\right](\mu)$ which is $\mathscr{C}^{\infty}$ in $\mathscr{X}^{\prime}$ and analytic in $\mu$. Since the pole of $\mathscr{N}(\mu)^{-1}$ is of order 1 , see [CoDa4], and since by construction, the columns of $\mathfrak{W}\left[x^{\prime}\right](\mu)$ belong to the special space $\mathfrak{H}_{0}^{\lambda}$ of homogeneous functions, we have obtained

Lemma C.4.1 The leading singular function $\mathscr{X}^{\prime} \mapsto \Phi^{0}\left[\mathscr{X}^{\prime}\right]$ is $\mathscr{C}^{\infty}(\mathscr{E})$ with values in $\mathfrak{H}_{0}^{\lambda}$, which means that there exists $N$-component vectors $\boldsymbol{q}_{0}^{ \pm}\left[x^{\prime}\right](\alpha)$ meromorphic in $\alpha$ and $\mathscr{C}^{\infty}$ in $\mathscr{X}^{\prime}$ such that

$$
\begin{equation*}
\Phi^{0}\left[\mathscr{x}^{\prime}\right]=\int_{|\alpha|=\delta}(\alpha \zeta+\bar{\zeta})^{\lambda} \boldsymbol{q}_{0}^{+}\left[\mathscr{x}^{\prime}\right](\alpha) \mathrm{d} \alpha+\int_{|\alpha|=\delta}(\zeta+\alpha \bar{\zeta})^{\lambda} \boldsymbol{q}_{0}^{-}\left[\mathscr{x}^{\prime}\right](\alpha) \mathrm{d} \alpha \tag{C.4.2}
\end{equation*}
$$

The first shadow singularity $\Phi^{1}$ is given by

$$
\begin{equation*}
\Phi^{1}\left[\mathscr{x}^{\prime}\right]=-\frac{1}{2 i \pi} \int_{\gamma(\lambda)+1} r^{\mu}\left[\mathfrak{A}^{0}\left[\mathscr{x}^{\prime}\right](\mu)\right]^{-1} \mathfrak{M}\left(\mathbb{1}_{r \in[0,1]}\left(r^{\beta} \mathfrak{A}^{1} \Phi^{0}\right)\right)\left[\mathscr{x}^{\prime}\right](\mu) \mathrm{d} \mu \tag{C.4.3}
\end{equation*}
$$

The following lemmas give that the structure of $\mathfrak{A}^{1} \Phi^{0}$ is compatible with representations of the type (C.3.7).

Lemma C.4.2 Let $\lambda \in \mathbb{C}$. For any $j \in \mathbb{N}$, the operator $L^{j}$ acts from $\mathscr{C}^{\infty}\left(\mathscr{E}^{\prime}, \mathfrak{H}_{0}^{\lambda}\right)$ into $\mathscr{C}^{\infty}\left(\mathscr{E}, \mathfrak{H}_{0}^{\lambda+j-2 d}\right)$.

Proof. The operator $L^{j}$ is a linear combination with $\mathscr{C}^{\infty}(\mathscr{E})$ coefficients of terms of the form $y^{\beta} \partial_{\mathscr{x}^{\prime}}^{i} \partial_{y}^{\delta}$ with $|\delta|-|\beta|=2 d-j$. The derivative $\partial_{x^{\prime}}^{i}$, acts only on the coefficients depending on $x^{\prime}$ and do not change the angular structure, so we may discard it. We are left with $y^{\beta} \partial_{y}^{\delta}$, which we can write as a linear combination of terms

$$
\zeta^{\beta_{1}} \bar{\zeta}^{\beta_{2}} \partial_{\zeta}^{\delta_{1}} \partial_{\bar{\zeta}}^{\delta_{2}} \quad \text { with } \quad \delta_{1}+\delta_{2}-\beta_{1}-\beta_{2}=2 d-j .
$$

It is clear that it suffices to prove that for any $\delta_{1}, \delta_{2}, \beta_{1}$ and $\beta_{2}$ with $\delta_{1}+\delta_{2}-\beta_{1}-\beta_{2}=$ $2 d-j$, and for any function $q(\alpha)$ meromorphic in $\alpha$, there exists $q^{\prime}(\alpha)$ also meromorphic in $\alpha$ such that

$$
\zeta^{\beta_{1}} \bar{\zeta}^{\beta_{2}} \partial_{\zeta}^{\delta_{1}} \partial_{\bar{\zeta}}^{\delta_{2}} \int_{|\alpha|=\delta}(\alpha \zeta+\bar{\zeta})^{\lambda} q(\alpha) \mathrm{d} \alpha=\int_{|\alpha|=\delta}(\alpha \zeta+\bar{\zeta})^{\lambda+j-2 d} q^{\prime}(\alpha) \mathrm{d} \alpha
$$

We have

$$
\partial_{\zeta}^{\delta_{1}} \partial_{\bar{\zeta}}^{\delta_{2}} \int_{|\alpha|=\delta}(\alpha \zeta+\bar{\zeta})^{\lambda} q(\alpha) \mathrm{d} \alpha=c \int_{|\alpha|=\delta} \alpha^{\delta_{1}}(\alpha \zeta+\bar{\zeta})^{\lambda-|\delta|} q(\alpha) \mathrm{d} \alpha
$$

With the equality $\bar{\zeta}=(\alpha \zeta+\bar{\zeta})-\alpha \zeta$, we transform $\zeta^{\beta_{1}} \bar{\zeta}^{\beta_{2}}$ into a linear combination of terms of the form $\zeta^{\gamma_{1}}(\alpha \zeta+\bar{\zeta})^{\gamma_{2}}$. Thus we are left with integrals of the form

$$
\int_{|\alpha|=\delta} \zeta^{n}(\alpha \zeta+\bar{\zeta})^{\lambda+j-2 d-n} q(\alpha) \mathrm{d} \alpha
$$

As $\partial_{\alpha}^{n}(\alpha \zeta+\bar{\zeta})^{\lambda+j-2 d}=c \zeta^{n}(\alpha \zeta+\bar{\zeta})^{\lambda+j-2 d-n}$, we integrate by parts $n$ times in the above integral and obtain the result.

In the same way, we obtain the corresponding result for the trace operators:

Lemma C.4.3 Let $\lambda \in \mathbb{C}$. For any $j \in \mathbb{N}_{0}$ the operator $B^{j}$ acts from $\mathscr{C}^{\infty}\left(\mathscr{E}, \mathfrak{H}_{0}^{\lambda}\right)$ into $\mathscr{C}^{\infty}\left(\mathscr{E}, \mathfrak{H}_{0}^{\lambda+j-\rho}\right)$, where $\mathfrak{H}_{0}^{\lambda-\rho}$ is the space of $m$-component functions homogeneous of degree $\left(\lambda-\rho_{1}, \ldots, \lambda-\rho_{m}\right)$ with Cayley representation like (C.3.7).

Let us return to (C.4.3). Let $\left(F^{1}, \boldsymbol{G}_{ \pm}^{1}\right)\left[\mathscr{X}^{\prime}\right]$ be the traces on $r=1$ of $\mathfrak{A}^{1}\left[\mathscr{X}^{\prime}\right] \Phi^{0}\left[\mathscr{X}^{\prime}\right]$. We have

$$
\mathfrak{M}\left(\mathbb{1}_{r \in[0,1]}\left(r^{\boldsymbol{\beta}} \mathfrak{A}^{1} \Phi^{0}\right)\right)\left[\mathscr{X}^{\prime}\right](\mu)=\frac{1}{\mu-(\lambda+1)}\left(\boldsymbol{F}^{1}, \boldsymbol{G}_{ \pm}^{1}\right)\left[\mathscr{X}^{\prime}\right] .
$$

By Lemma C.4.2, $r^{\lambda+1-2 d} \boldsymbol{F}^{1}\left[\mathscr{X}^{\prime}\right](\theta)$ belongs to $\mathscr{C}^{\infty}\left(\mathscr{E}, \mathfrak{H}_{0}^{\lambda+1-2 d}\right)$ : There exists $\boldsymbol{q}^{ \pm}\left[\mathscr{X}^{\prime}\right]$ such that
$r^{\lambda+1-2 d} \boldsymbol{F}^{1}\left[\mathscr{X}^{\prime}\right]=\int_{|\alpha|=\delta}(\alpha \zeta+\bar{\zeta})^{\lambda+1-2 d} \boldsymbol{q}^{+}\left[\mathscr{X}^{\prime}\right](\alpha) \mathrm{d} \alpha+\int_{|\alpha|=\delta}(\zeta+\alpha \bar{\zeta})^{\lambda+1-2 d} \boldsymbol{q}^{-}\left[\mathscr{X}^{\prime}\right](\alpha) \mathrm{d} \alpha$. We define for $\mu \in \mathbb{C}$ the following element $\boldsymbol{f}_{0}^{1}\left[x^{\prime}\right](\mu) \in \mathfrak{H}_{0}^{\mu-2 d}$ :

$$
\boldsymbol{f}_{0}^{1}\left[\mathscr{X}^{\prime}\right](\mu):=\int_{|\alpha|=\delta}(\alpha \zeta+\bar{\zeta})^{\mu-2 d} \boldsymbol{q}^{+}\left[\mathscr{X}^{\prime}\right](\alpha) \mathrm{d} \alpha+\int_{|\alpha|=\delta}(\zeta+\alpha \bar{\zeta})^{\mu-2 d} \boldsymbol{q}^{-}\left[\mathscr{X}^{\prime}\right](\alpha) \mathrm{d} \alpha
$$

Of course, $\boldsymbol{f}_{0}^{1}\left[\mathscr{X}^{\prime}\right](\lambda+1)=r^{\lambda+1-2 d} \boldsymbol{F}^{1}\left[\mathscr{X}^{\prime}\right]$. Let us denote

$$
\boldsymbol{f}^{1}\left[\mathscr{X}^{\prime}\right](\mu):=r^{\mu-2 d} \boldsymbol{F}^{1}\left[x^{\prime}\right]-\boldsymbol{f}_{0}^{1}\left[x^{\prime}\right](\mu) .
$$

It is clear that in the representation formula (C.3.6), we may take as right inverse of $r^{\mu-2 d} \boldsymbol{F}^{1}$,

$$
\mathfrak{R}_{0}\left[\mathscr{X}^{\prime}\right](\mu)\left(\boldsymbol{f}_{0}^{1}\left[\mathscr{X}^{\prime}\right](\mu)\right)+\mathfrak{R}\left[\mathscr{X}^{\prime}\right](\mu)\left(\boldsymbol{f}^{1}\left[\mathscr{x}^{\prime}\right](\mu)\right),
$$

instead of $\mathfrak{R}\left[\mathscr{X}^{\prime}\right](\mu)\left(r^{\mu-2 d} \boldsymbol{F}^{1}\right)$. Therefore we have the following decomposition in four parts of the integrand of (C.4.3):

$$
\begin{align*}
r^{\mu}\left[\mathfrak{A}^{0}\left[\mathscr{X}^{\prime}\right](\mu)\right]^{-1} \mathfrak{M}\left(\mathbb{1}_{r \in[0,1]}\left(r^{\boldsymbol{\beta}} \mathfrak{A}^{1} \Phi^{0}\right)\right)(\mu) & =\frac{1}{\mu-(\lambda+1)}\left(r^{\mu}\left[\mathfrak{A}^{0}\left[\mathscr{x}^{\prime}\right](\mu)\right]^{-1}\left(\boldsymbol{F}^{1}, \boldsymbol{G}_{ \pm}^{1}\right)\right) \\
& =\frac{1}{\mu-(\lambda+1)}\left(U_{1}+U_{2}+U_{3}+U_{4}\right)(\mu) . \tag{C.4.4}
\end{align*}
$$

where

$$
\begin{aligned}
U_{1}(\mu) & =\mathfrak{R}\left[\mathscr{x}^{\prime}\right](\mu) \boldsymbol{f}^{1}\left[x^{\prime}\right](\mu), \\
U_{2}(\mu) & =\mathfrak{R}_{0}\left[\mathscr{x}^{\prime}\right](\mu) \boldsymbol{f}_{0}^{1}\left[\mathscr{x}^{\prime}\right](\mu), \\
U_{3}(\mu) & =\mathfrak{W}\left[\mathscr{x}^{\prime}\right](\mu) \mathscr{N}\left[\mathscr{X}^{\prime}\right](\mu)^{-1}\left(-\mathfrak{g}_{ \pm} B^{0}\left[\mathscr{x}^{\prime}\right] \mathfrak{R}\left[\mathscr{x}^{\prime}\right](\mu) \boldsymbol{f}^{1}\left[\mathscr{X}^{\prime}\right](\mu)\right), \\
U_{4}(\mu) & =\mathfrak{W}\left[\mathscr{X}^{\prime}\right](\mu) \mathscr{N}\left[\mathscr{x}^{\prime}\right](\mu)^{-1}\left(\boldsymbol{G}_{ \pm}^{1}-\mathfrak{g}_{ \pm} B^{0}\left[\mathscr{X}^{\prime}\right] \mathfrak{R}_{0}\left[\mathscr{x}^{\prime}\right](\mu) \boldsymbol{f}_{0}^{1}\left[\mathscr{X}^{\prime}\right](\mu)\right) .
\end{aligned}
$$

Coming back to (C.4.3), we have to compute the contour integral

$$
\Phi^{1}\left[\mathscr{X}^{\prime}\right]=-\frac{1}{2 i \pi} \int_{\gamma(\lambda)+1} \frac{1}{\mu-(\lambda+1)}\left(U_{1}+U_{2}+U_{3}+U_{4}\right)(\mu) \mathrm{d} \mu .
$$

Let us compute the residue in $\mu=\lambda+1$ of each of the four terms.
(i) As $\boldsymbol{f}^{1}\left[\mathscr{X}^{\prime}\right](\lambda+1)=0$, the residue of $(\mu-(\lambda+1))^{-1} U_{1}(\mu)$ is 0 .
(ii) The residue of $(\mu-(\lambda+1))^{-1} U_{2}(\mu)$ is equal to $U_{2}(\lambda+1)$, which coincides with $\mathfrak{R}_{0}\left[\mathscr{X}^{\prime}\right](\lambda+1) \boldsymbol{f}_{0}^{1}\left[\mathscr{X}^{\prime}\right](\lambda+1)$, therefore belongs to $\mathfrak{H}_{0}^{\lambda+1}$.
(iii) As $\boldsymbol{f}^{1}\left[\mathscr{X}^{\prime}\right](\lambda+1)=0$, the pole of $(\mu-(\lambda+1))^{-1} U_{3}(\mu)$ is of order 1 , and the residue is a linear combination of the $\boldsymbol{w}_{\ell}^{ \pm}\left[\mathscr{x}^{\prime}\right](\lambda+1)$, therefore belongs to $\mathfrak{H}_{0}^{\lambda+1}$.
(iv) Finally, the pole of $(\mu-(\lambda+1))^{-1} U_{4}(\mu)$ in $\lambda+1$ is, a priori, of order 2:

$$
\begin{equation*}
\frac{1}{2 i \pi} \int_{\gamma(\lambda)+1} \mathfrak{W}\left[\mathscr{X}^{\prime}\right](\mu) \frac{\mathscr{N}\left[\mathscr{X}^{\prime}\right](\mu)^{-1}}{\mu-(\lambda+1)}\left(G_{ \pm}^{1}-\mathfrak{g}_{ \pm} B^{0}\left[\mathscr{X}^{\prime}\right] \mathfrak{R}_{0}\left[\mathscr{X}^{\prime}\right](\mu) \boldsymbol{f}_{0}^{1}\left[\mathscr{X}^{\prime}\right](\mu)\right) \mathrm{d} \mu \tag{C.4.5}
\end{equation*}
$$

The term (C.4.5) is itself the sum of an element of $\mathfrak{H}_{0}^{\lambda+1}$, cf. (iii), and of

$$
\begin{equation*}
\frac{1}{2 i \pi} \int_{\gamma(\lambda)+1} \mathfrak{W}\left[\mathscr{x}^{\prime}\right](\mu) \frac{\mathscr{Y}\left[\mathscr{x}^{\prime}\right](\mu)^{-1}}{\mu-(\lambda+1)}\left(G_{ \pm}^{1}-\mathfrak{g}_{ \pm} B^{0}\left[\mathscr{X}^{\prime}\right] \mathfrak{R}_{0}\left[\mathscr{X}^{\prime}\right](\lambda+1)\left(r^{\lambda+1-2 d} \boldsymbol{F}^{1}\right)\right) \mathrm{d} \mu \tag{C.4.6}
\end{equation*}
$$

By construction $G_{ \pm}^{1}$ is the couple of traces $\mathfrak{g}_{ \pm} B^{1}\left[\mathscr{X}^{\prime}\right] \Phi^{0}\left[\mathscr{X}^{\prime}\right]$. Therefore

$$
\boldsymbol{G}_{ \pm}^{1}-\mathfrak{g}_{ \pm} B^{0}\left[\mathscr{x}^{\prime}\right] \mathfrak{R}_{0}\left[\mathscr{x}^{\prime}\right](\lambda+1)\left(r^{\lambda+1-2 d} \boldsymbol{F}\right)=\mathfrak{g}_{ \pm} \Psi^{1}\left[\mathscr{x}^{\prime}\right]
$$

where

$$
\Psi^{1}\left[\mathscr{x}^{\prime}\right]:=B^{1}\left[\mathscr{x}^{\prime}\right] \Phi^{0}\left[\mathscr{X}^{\prime}\right]-B^{0}\left[x^{\prime}\right] \mathfrak{R}_{0}\left[x^{\prime}\right](\lambda+1)\left(L^{1}\left[\mathscr{x}^{\prime}\right] \Phi^{0}\left[\mathscr{x}^{\prime}\right]\right) .
$$

The $m$-component function $\Psi^{1}\left[\mathscr{X}^{\prime}\right]$ belongs to $\mathscr{C}^{\infty}\left(\mathscr{E}, \mathfrak{H}_{0}^{\lambda+1-\rho}\right)$ by virtue of Lemmas C.4.2 and C.4.3. Gathering the results for $\Phi^{1}$, we have obtained

Lemma C.4.4 The first shadow singularity $\Phi^{1}\left[x^{\prime}\right]$ is the sum of $\Phi_{0}^{1}\left[\mathscr{X}^{\prime}\right]$ which belongs to $\mathscr{C}^{\infty}\left(\mathscr{E}, \mathfrak{H}_{0}^{\lambda+1}\right)$ and of $\Phi_{1}^{1}$ :

$$
\begin{equation*}
\Phi_{1}^{1}:=\frac{1}{2 i \pi} \int_{\gamma(\lambda)+1} \mathfrak{W}\left[\mathscr{x}^{\prime}\right](\mu) \frac{\mathscr{N}\left[\mathscr{x}^{\prime}\right](\mu)^{-1}}{\mu-(\lambda+1)}\left(\mathfrak{g}_{ \pm} \Psi^{1}\left[\mathscr{x}^{\prime}\right]\right) \mathrm{d} \mu \tag{С.4.7}
\end{equation*}
$$

where $\Psi^{1}\left[\mathscr{X}^{\prime}\right]$ belongs to $\mathscr{C}^{\infty}\left(\mathscr{E}, \mathfrak{H}_{0}^{\lambda+1-\rho}\right)$.

## C. 5 The relation of compatibility

Our aim is to show that the coefficient in front of the term $(\mu-(\lambda+1))^{-2}$ in the Laurent expansion of $(\mu-(\lambda+1))^{-1} \mathscr{N}\left[\mathscr{X}^{\prime}\right](\mu)^{-1}\left(\mathfrak{g}_{ \pm} \Psi^{1}\left[\mathscr{X}^{\prime}\right]\right)$ is zero. As $\mathscr{N}\left[\mathscr{X}^{\prime}\right](\mu)^{-1}$ has its a pole of order 1 in $\lambda+1$, the necessary and sufficient condition for this coefficient to be zero is that

$$
\begin{equation*}
\mathfrak{g}_{ \pm} \Psi^{1}\left[\mathscr{X}^{\prime}\right] \in \operatorname{rg} \mathscr{N}\left[\mathscr{X}^{\prime}\right](\lambda+1), \tag{C.5.1}
\end{equation*}
$$

which is the "relation of compatibility".

Lemma C.5. 1 Let $\lambda$ be of the form $\frac{1}{2}+k$ with $k \in \mathbb{Z}$. Let $x^{\prime} \in \mathscr{E}$. Then the range of $\mathscr{N}\left[\mathscr{X}^{\prime}\right](\lambda)$ is the subspace of the $\left(b_{-}^{1}, \ldots, b_{-}^{m}, b_{+}^{1}, \ldots, b_{+}^{m}\right)$ which satisfy $b_{-}^{h}=-b_{+}^{h}$ for $h=1, \ldots, m$.

Proof. Let us fix $x^{\prime}$ and let us drop it in the notations. In the case when the roots $\alpha_{\ell}^{ \pm}$are distinct, according to [CoDa4, §3], $\mathscr{N}(\mu)$ has the general structure, by $m \times m$ blocks:

$$
\mathscr{N}(\mu)=\left(\begin{array}{cc}
E(\mu) & 0 \\
0 & E(\mu)
\end{array}\right)\left(\begin{array}{cc}
e^{-i \pi \mu} \mathfrak{B}^{+} & -e^{-i \pi \mu} \mathfrak{B}^{-} \\
-e^{-i \pi \mu} \mathfrak{B}^{+} & e^{-i \pi \mu} \mathfrak{B}^{-}
\end{array}\right)\left(\begin{array}{cc}
F^{+}(\mu) & 0 \\
0 & F^{-}(\mu)
\end{array}\right),
$$

where $E(\mu)$ is a diagonal matrix everywhere invertible except on a finite number of integers, $F^{ \pm}(\mu)$ are everywhere invertible and the two matrices $\mathfrak{B}^{ \pm}$are invertible, due to the ellipticity of the boundary value problem, see [CoDa4, $\S 4]$. The statement of the lemma for $\mu=\lambda$ is straightforward in this case. The general case where the $\alpha_{\ell}^{ \pm}$are not supposed distinct is obtained by perturbation.

Lemma C.5.2 Let $\lambda$ be of the form $\frac{1}{2}+k$ with $k \in \mathbb{Z}$. Let $\Psi$ belong to $\mathfrak{H}_{0}^{\lambda-\rho}$. Then $\mathfrak{g}_{-} \Psi=-\mathfrak{g}_{+} \Psi$.

Proof. Let $\Psi_{h}$ denote the components of $\Psi$, for $h=1, \ldots, m$. The component $\Psi_{h}$ belongs to $\mathfrak{H}_{0}^{\lambda-\rho_{h}}$, which means that there exists functions $p_{h}^{ \pm}$meromorphic in $\alpha$ and such that

$$
\Psi_{h}=\int_{|\alpha|=\delta}(\alpha \zeta+\bar{\zeta})^{\lambda-\rho_{h}} p_{h}^{+}(\alpha) \mathrm{d} \alpha+\int_{|\alpha|=\delta}(\zeta+\alpha \bar{\zeta})^{\lambda-\rho_{h}} p_{h}^{-}(\alpha) \mathrm{d} \alpha
$$

It remains to compute the traces $\mathfrak{g}_{ \pm}$of $\Psi_{h}$. We use the formulae

$$
(\alpha \zeta+\bar{\zeta})^{\mu}=\bar{\zeta}^{\mu}\left(1+\alpha \frac{\zeta}{\bar{\zeta}}\right)^{\mu} \quad \text { and } \quad(\zeta+\alpha \bar{\zeta})^{\mu}=\zeta^{\mu}\left(1+\alpha \frac{\bar{\zeta}}{\zeta}\right)^{\mu}
$$

There holds (since $|\alpha|<1$ )

$$
\zeta^{\mu}=r e^{i \mu \theta}, \quad \bar{\zeta}^{\mu}=r e^{-i \mu \theta}, \quad\left(1+\alpha \frac{\zeta}{\bar{\zeta}}\right)^{\mu}=\left(1+\alpha e^{2 i \theta}\right)^{\mu}, \quad\left(1+\alpha \frac{\bar{\zeta}}{\bar{\zeta}}\right)^{\mu}=\left(1+\alpha e^{-2 i \theta}\right)^{\mu}
$$

Whence

$$
\begin{aligned}
& \mathfrak{g}_{-} \Psi_{h}=e^{i\left(\lambda-\rho_{h}\right) \pi} \int_{|\alpha|=\delta}(1+\alpha)^{\lambda-\rho_{h}} p_{h}^{+}(\alpha) \mathrm{d} \alpha+e^{-i\left(\lambda-\rho_{h}\right) \pi} \int_{|\alpha|=\delta}(1+\alpha)^{\lambda-\rho_{h}} p_{h}^{-}(\alpha) \mathrm{d} \alpha \\
& \mathfrak{g}_{+} \Psi_{h}=e^{-i\left(\lambda-\rho_{h}\right) \pi} \int_{|\alpha|=\delta}(1+\alpha)^{\lambda-\rho_{h}} p_{h}^{+}(\alpha) \mathrm{d} \alpha+e^{i\left(\lambda-\rho_{h}\right) \pi} \int_{|\alpha|=\delta}(1+\alpha)^{\lambda-\rho_{h}} p_{h}^{-}(\alpha) \mathrm{d} \alpha
\end{aligned}
$$

As $\lambda=\frac{1}{2}+k$, we have obtained the lemma.
The consequence of Lemmas C.5.1 and C.5.2 for $\Phi^{1}\left[x^{\prime}\right]$ is now clear: (C.5.1) holds. Therefore the function $\Psi_{1}^{1}\left[\mathscr{X}^{\prime}\right]$ defined in (C.4.7) also belongs to $\mathscr{C}^{\infty}\left(\mathscr{E}, \mathfrak{H}_{0}^{\lambda+1}\right)$. Which means that, finally, the first shadow singularity $\Phi^{1}\left[\mathscr{X}^{\prime}\right]$ belongs to $\mathscr{C}^{\infty}\left(\mathscr{E}, \mathfrak{H}_{0}^{\lambda+1}\right)$, i.e. satisfies at its degree of homogeneity exactly the same property as $\Phi^{0}\left[\mathscr{X}^{\prime}\right]$, see Lemma C.4.1.

The proof of this property can be immediately generalized to the following:
Proposition C.5.3 Let $\lambda \in \mathbb{C}$ of the form $\frac{1}{2}+k$ with integer $k$. Let $F\left[x^{\prime}\right](r, \theta)$ belong to $\mathscr{C}^{\infty}\left(\mathscr{E}, \mathfrak{H}_{0}^{\lambda-2 d}\right)$ and $G_{ \pm}\left[\mathscr{X}^{\prime}\right](r)$ be the traces on $\theta= \pm \pi$ of a $m$-component vector function $\Psi\left[x^{\prime}\right] \in \mathscr{C}^{\infty}\left(\mathscr{E}, \mathfrak{H}_{0}^{\lambda-\rho}\right)$. Then the $N$-component function $\Phi\left[x^{\prime}\right]$ defined as

$$
\Phi\left[\mathscr{X}^{\prime}\right]=\frac{1}{2 i \pi} \int_{\gamma(\lambda)} r^{\mu}\left[\mathfrak{A}^{0}\left[\mathscr{x}^{\prime}\right](\mu)\right]^{-1} \mathfrak{M}\left(\mathbb{1}_{r \in[0,1]}\left(r^{2 d} F, r^{\boldsymbol{\rho}} G\right)\right)\left[\mathscr{X}^{\prime}\right](\mu) \mathrm{d} \mu
$$

belongs to $\mathscr{C}^{\infty}\left(\mathscr{E}, \mathfrak{H}_{0}^{\lambda}\right)$.
Therefore, with the help of Lemmas C.4.2 and C.4.3, we see that the procedure for the analysis of the successive shadows $\Phi^{2}, \ldots, \Phi^{p}$ is recursive. Therefore for all $p \in \mathbb{N}_{0}, \Phi^{p}$ belong to $\mathscr{C}^{\infty}\left(\mathscr{E}, \mathfrak{H}_{0}^{\lambda+p}\right)$ and, thus, do not contain any logarithmic term.

## C. 6 Absence of logarithms, general results

Examining the arguments of the proofs of Lemmas C.4.1 to C.4.3 and Proposition C.5.3, we can see that, in fact, the result we have proved does not use any ellipticity in the edge variable $\mathscr{X}^{\prime} \in \mathscr{E}$, only the smooth dependency. In the next statement, we select the hypotheses which are sufficient to obtain our result on the absence of logarithms in shadow singularities: Hypothesis C.6.1 Let $\mathscr{X}^{\prime} \mapsto\left(L^{0}, B^{0}\right)\left[\mathscr{X}^{\prime}\right]$ be $\mathscr{C}^{\infty}(\mathscr{E})$ with values in the space $\mathrm{Op}_{\mathrm{Ell}}^{2 d, \boldsymbol{\rho}}\left(\mathbb{R}^{2}\right)$ of $(N \times N)$ elliptic systems homogeneous of order $2 d$ with constant coefficients in $\mathbb{R}^{2}$, with complementing boundary conditions homogeneous of degree $\boldsymbol{\rho}=\left(\rho_{1}, \ldots, \rho_{m}\right)$ with constant coefficients. The Mellin symbol of $\left(L^{0}, \gamma_{ \pm} B^{0}\right)\left[x^{\prime}\right]$ is denoted by $\mathfrak{A}^{0}\left[x^{\prime}\right]$ with $\gamma_{-}$ and $\gamma_{+}$the traces on $\left\{\left(y_{1}, y_{2}\right) \mid y_{1}<0\right\}$ from below and from above respectively.

For any $j \in \mathbb{N}$, let $x^{\prime} \mapsto\left(L^{j}, B^{j}\right)\left[x^{\prime}\right]$ be a matrix-function with coefficients $L_{k, \ell}^{j}\left[\mathscr{X}^{\prime}\right]$ and $B_{h, \ell}^{j}\left[x^{\prime}\right], \mathscr{C}^{\infty}(\mathscr{E})$ with values in the space of operators

$$
\mathrm{Op}^{2 d-j}\left(\mathbb{R}^{2}\right) \text { for } L_{k, \ell}^{j} \text { and } \mathrm{Op}^{\rho_{h}-j}\left(\mathbb{R}^{2}\right) \text { for } B_{h, \ell}^{j},
$$

where for $p \in \mathbb{Z}, \mathrm{Op}^{p}\left(\mathbb{R}^{2}\right)$ is defined as the space of finite linear combinations with $\mathscr{C}^{\infty}(\mathscr{E})$ coefficients of partial differential operators of the form $y^{\beta} \partial_{\mathscr{x}}^{i}, \partial_{y}^{\delta}$ with $|\delta|-|\beta|=p$. We denote the triple $\left(L^{j}, \gamma_{ \pm} B^{j}\right)\left[\mathscr{X}^{\prime}\right]$ by $\mathfrak{A l}^{j}\left[\mathscr{x}^{\prime}\right]$.

The proofs of Lemmas C.4.1 to C.4.3 and Proposition C.5.3 then yield
Theorem C.6.2 Let $\left(L^{j}, B^{j}\right)_{j \geq 0}$ be a sequence of operators satisfying Hypothesis C.6.1. Let $\lambda=\frac{1}{2}+k$ with $k \in \mathbb{Z}$ and let $\gamma(\lambda)$ be the circle with center $\lambda$ and radius $\frac{1}{4}$. With the function $\left(\mathscr{X}^{\prime}, \mu\right) \mapsto\left(F, G_{ \pm}\right)\left[\mathscr{X}^{\prime}\right](\mu)$ supposed to be $\mathscr{C}^{\infty}(\mathscr{E})$ in $\mathscr{X}^{\prime}$ and analytic in $\mu$, with values in $\mathrm{L}^{2}(-\pi, \pi) \times \mathbb{C}^{m} \times \mathbb{C}^{m}$, we define the following leading singularity, which is a generalization of (C.2.3):

$$
\Phi^{0}\left[x^{\prime}\right]=\frac{1}{2 i \pi} \int_{\gamma(\lambda)} r^{\mu}\left[\mathfrak{A}^{0}\left[\mathscr{x}^{\prime}\right](\mu)\right]^{-1}\left(\boldsymbol{F}, \boldsymbol{G}_{ \pm}\right)\left[\mathscr{x}^{\prime}\right](\mu) \mathrm{d} \mu
$$

and its sequence of shadows $\left(\Phi^{p}\left[\mathscr{x}^{\prime}\right]\right)_{p}$ according to Definition C.2.1. Then, for any integer $p \geq 0, \Phi^{p}\left[\mathscr{x}^{\prime}\right]$ belongs to $\mathscr{C}^{\infty}\left(\mathscr{E}, \mathfrak{H}_{0}^{\lambda+p}\right)$. In particular $\Phi^{p}\left[\mathscr{x}^{\prime}\right](r, \theta)$ can be written in the form $r^{\lambda+p} \boldsymbol{\psi}\left(x^{\prime}, \theta\right)$ with $\boldsymbol{\psi} \in \mathscr{C}^{\infty}(\mathscr{E} \times[-\pi, \pi]) \otimes \mathbb{C}^{N}$.

In fact this statement extends to the wider class of Agmon-Douglis-Nirenberg systems with covering boundary conditions:
Hypothesis C.6.3 Let $N \in \mathbb{N}, \boldsymbol{\sigma}=\left(\sigma_{1}, \ldots, \sigma_{N}\right), \boldsymbol{\tau}=\left(\tau_{1}, \ldots, \tau_{N}\right)$,

$$
m=\frac{1}{2}\left(\sigma_{1}-\tau_{1}+\ldots+\sigma_{N}-\tau_{N}\right)
$$

and $\boldsymbol{\rho}=\left(\rho_{1}, \ldots, \rho_{m}\right)$. Let $x^{\prime} \mapsto\left(L^{0}, B^{0}\right)\left[x^{\prime}\right]$ be $\mathscr{C}^{\infty}(\mathscr{E})$ with values in the space $\mathrm{Op}_{\mathrm{ADN}}^{\boldsymbol{\sigma}, \boldsymbol{\tau}, \boldsymbol{\rho}}\left(\mathbb{R}^{2}\right)$ of $(N \times N)$ Agmon-Douglis-Nirenberg elliptic systems homogeneous of or$\operatorname{der} \sigma_{k}-\tau_{\ell}$ with constant coefficients in $\mathbb{R}^{2}$, with complementing boundary conditions homogeneous of degree $\rho_{h}-\tau_{\ell}$ with constant coefficients.

For any $j \in \mathbb{N}$, let $\mathscr{X}^{\prime} \mapsto\left(L^{j}, B^{j}\right)\left[\mathscr{X}^{\prime}\right]=: \mathfrak{A}^{j}\left[\mathscr{X}^{\prime}\right]$ be a matrix-function with coefficients $L_{k, \ell}^{j}\left[\mathscr{X}^{\prime}\right]$ and $B_{h, \ell}^{j}\left[\mathscr{x}^{\prime}\right], \mathscr{C}^{\infty}(\mathscr{E})$ with values in the space of operators

$$
\mathrm{Op}^{\sigma_{k}-\tau_{\ell}-j}\left(\mathbb{R}^{2}\right) \text { for } L_{k, \ell}^{j} \text { and } \mathrm{Op}^{\rho_{h}-\tau_{\ell}-j}\left(\mathbb{R}^{2}\right) \text { for } B_{h, \ell}^{j},
$$

with $\mathrm{Op}^{p}\left(\mathbb{R}^{2}\right)$ as in Hypothesis C.6.1.
The Mellin transform and the Cayley representation can be used with the same success in the framework of Agmon-Douglis-Nirenberg systems, see [CoDa2, CoDa4], which allows to obtain:

Theorem C.6.4 Let $\left(L^{j}, B^{j}\right)_{j \geq 0}$ be a sequence of operators satisfying Hypothesis C.6.3. Let $\lambda=\frac{1}{2}+k, \gamma(\lambda)$ and $\left(\boldsymbol{F}, \boldsymbol{G}_{ \pm}\right)\left[\mathscr{x}^{\prime}\right](\mu)$ be as in Theorem C.6.2. We define the following leading singularity:

$$
\Phi^{0}\left[\mathscr{X}^{\prime}\right]=\frac{1}{2 i \pi} \int_{\gamma(\lambda)} r^{\mu-\boldsymbol{\tau}}\left[\mathfrak{A}^{0}\left[\mathscr{X}^{\prime}\right](\mu)\right]^{-1}\left(\boldsymbol{F}, \boldsymbol{G}_{ \pm}\right)\left[\mathscr{X}^{\prime}\right](\mu) \mathrm{d} \mu,
$$

and its sequence of shadows $\left(\Phi^{p}\left[\mathscr{X}^{\prime}\right]\right)_{p}$ by an obvious modification of Definition C.2.1, with $\boldsymbol{\beta}=\left(\sigma_{1}, \ldots, \sigma_{N}, \rho_{1}, \ldots, \rho_{m}, \rho_{1}, \ldots, \rho_{m}\right)$ and $\mu$ replaced with $\mu-\boldsymbol{\tau}$ as above.

Then $\Phi^{p}\left[\mathscr{x}^{\prime}\right]$ is homogeneous of multi-degree $\lambda+p-\boldsymbol{\tau}$, i.e. its $j$-th component $\Phi_{j}^{p}$ satisfies $\Phi_{j}^{p}\left[\mathscr{X}^{\prime}\right](r, \theta)=r^{\lambda+p-\tau_{j}} \psi_{j}\left(\mathscr{X}^{\prime}, \theta\right)$ with $\psi_{j} \in \mathscr{C}^{\infty}(\mathscr{E} \times[-\pi, \pi])$.

We obtain as a corollary (and a generalization of Theorem A.4.3) that the asymptotics along a crack edge of the solutions of Agmon-Douglis-Nirenberg systems associated with coercive bilinear forms contain no logarithmic term:

Corollary C.6.5 Let $(L, B)$ be an $(N \times N)$ Agmon-Douglis-Nirenberg elliptic system of order $\sigma_{k}-\tau_{\ell}$ with smooth coefficients in $\mathbb{R}^{n+1}$, with complementing boundary conditions homogeneous of degree $\rho_{h}-\tau_{\ell}$ with smooth coefficients. Let us assume that $(L, B)$ is associated with a coercive bilinear form. Let $\rho_{\max }:=\max \left\{\rho_{1}, \ldots, \rho_{m}\right\}$. Any solution $\boldsymbol{u}$ of problem (C.2.1) (with a smooth right hand side $\boldsymbol{f}$ ) which belongs to $\mathrm{H}^{s-\tau}\left(W_{\pi}\right)$ with
$s \geq \rho_{\max }$ has the following asymptotic expansion as $r \rightarrow 0$ : For any integer $K>k_{0}$

$$
\begin{array}{r}
\boldsymbol{u}=\sum_{j=1}^{m} c_{j}^{0}\left(\mathscr{X}^{\prime}\right) r^{\frac{1}{2}+k_{0}-\boldsymbol{\tau}} \boldsymbol{\psi}_{j}^{0}\left(\mathscr{X}^{\prime}, \theta\right)+\sum_{k=k_{0}+1}^{K} \sum_{j=1}^{j(k)} c_{j}^{k}\left(\mathscr{X}^{\prime}\right) r^{\frac{1}{2}+k-\boldsymbol{\tau}} \boldsymbol{\psi}_{j}^{k}\left(\mathscr{X}^{\prime}, \theta\right)  \tag{C.6.1}\\
+\boldsymbol{u}_{\mathrm{reg}, K}+\boldsymbol{u}_{\mathrm{rem}, K}
\end{array}
$$

where $k_{0}$ is the smallest integer such that $\frac{1}{2}+k_{0}>s-1$. The regular part $\boldsymbol{u}_{\mathrm{reg}, K}$ is in $\mathscr{C}^{\infty}\left(\mathbb{R}^{n+1}\right)$. The remainder $\boldsymbol{u}_{\mathrm{rem}, K}$ belongs to $\mathscr{C}^{K+1-\tau}\left(\bar{W}_{\pi}\right)$ and is flat of order $K-\boldsymbol{\tau}$ near $\mathscr{E}$.

## C. 7 ANGULAR DESCRIPTION OF SINGULAR FUNCTIONS

For simplicity, let us go back to the situation where Hypothesis C.6.1 is satisfied and let us consider $\Phi^{0}\left[\mathscr{X}^{\prime}\right]$ like in Theorem C.6.2, as well as its sequence of shadows $\left(\Phi^{p}\left[\mathscr{X}^{\prime}\right]\right)_{p}$. Theorem C.6.2 tells us that $\Phi^{p}\left[\mathscr{X}^{\prime}\right]$ belongs to $\mathscr{C}^{\infty}\left(\mathscr{E}, \mathfrak{H}_{0}^{\lambda+p}\right)$, which means that there exist meromorphic $\alpha \mapsto \boldsymbol{q}^{ \pm}\left[\mathscr{X}^{\prime}\right](\alpha)$ (with $\mathscr{C}^{\infty}(\mathscr{E})$ dependence on $\mathscr{X}^{\prime}$ ) such that

$$
\Phi^{p}\left[\mathscr{x}^{\prime}\right]=\int_{|\alpha|=\delta}(\alpha \zeta+\bar{\zeta})^{\lambda+p} \boldsymbol{q}^{+}\left[\mathscr{X}^{\prime}\right](\alpha) \mathrm{d} \alpha+\int_{|\alpha|=\delta}(\zeta+\alpha \bar{\zeta})^{\lambda+p} \boldsymbol{q}^{-}\left[\mathscr{x}^{\prime}\right](\alpha) \mathrm{d} \alpha
$$

But, in fact, the vector-functions $\boldsymbol{q}^{ \pm}\left[\mathscr{x}^{\prime}\right]$ are not arbitrary meromorphic functions in the unit disc: their poles belong to the set of the roots $\left\{\alpha_{\ell}^{ \pm}\left[\mathscr{X}^{\prime}\right]\right\}_{\ell=1, \ldots, m_{ \pm}}, c f$ (C.3.4).

As a consequence, as we are going to show, it is possible to give a modular representation of the $\Phi^{p}\left[x^{\prime}\right]$, if we assume

$$
\begin{equation*}
\text { The multiplicities } n_{\ell}^{ \pm} \text {of } \alpha_{\ell}^{ \pm}\left[x^{\prime}\right] \text { are constant on } \mathscr{E} \text {. } \tag{C1}
\end{equation*}
$$

Let $\boldsymbol{\alpha}\left[\mathscr{X}^{\prime}\right]$ denote the set $\left\{\left(\alpha_{\ell}^{ \pm}\left[\mathscr{X}^{\prime}\right], n_{\ell}^{ \pm}\right)\right\}$of the roots with their multiplicities.
Definition C.7.1 Under hypothesis $\left(\mathfrak{H}_{\mathrm{C} 1}\right)$, for any $\mu \in \mathbb{C}$ and $p \in \mathbb{N}$, let us denote by $\mathscr{C}^{\infty}\left(\mathscr{E}, \mathfrak{H}_{\boldsymbol{\alpha}, p}^{\mu}\right)$ the subspace of homogeneous functions $f\left[\mathscr{X}^{\prime}\right] \in \mathfrak{H}^{\mu}$ which admit a representation as:

$$
\begin{equation*}
f\left[\mathscr{X}^{\prime}\right]=\int_{|\alpha|=\delta}(\alpha \zeta+\bar{\zeta})^{\mu} q^{+}\left[\mathscr{X}^{\prime}\right](\alpha) \mathrm{d} \alpha+\int_{|\alpha|=\delta}(\zeta+\alpha \bar{\zeta})^{\mu} q^{-}\left[\mathscr{X}^{\prime}\right](\alpha) \mathrm{d} \alpha \tag{C.7.1}
\end{equation*}
$$

where the functions $q^{+}\left[\mathscr{X}^{\prime}\right]$ and $q^{-}\left[\mathscr{X}^{\prime}\right]$ are meromorphic in $\alpha, \mathscr{C}^{\infty}$ in $\mathscr{x}^{\prime}$, with poles only in the roots $\alpha_{\ell}^{+}\left[x^{\prime}\right]$ of order $\leq p n_{\ell}^{+}$and $\alpha_{\ell}^{-}\left[\mathscr{X}^{\prime}\right]$ of order $\leq p n_{\ell}^{-}$respectively. Let $\mathscr{C}^{\infty}\left(\mathscr{E}, \mathfrak{H}_{\alpha, p}^{\mu}\right)$ be the space $\mathscr{C}^{\infty}\left(\mathscr{E}, \mathfrak{H}_{\alpha, p}^{\mu}\right) \otimes \mathbb{C}^{N}$.

With these definitions, we have the following properties
(i) By (C.3.5), the kernel elements $\boldsymbol{w}_{\ell}^{ \pm}(\mu)$ belong to $\mathscr{C}^{\infty}\left(\mathscr{E}, \mathfrak{H}_{\alpha, 1}^{\mu}\right)$.
(ii) The definition (C.3.8) of $\mathfrak{R}_{0}$ gives us that for any $\mu \in \mathbb{C}$ and any $p \in \mathbb{N}$ $\mathfrak{R}_{0}$ acts from $\mathscr{C}^{\infty}\left(\mathscr{E}, \mathfrak{H}_{\alpha, p}^{\mu-2 d}\right)$ into $\mathscr{C}^{\infty}\left(\mathscr{E}, \mathfrak{H}_{\alpha, p+1}^{\mu}\right)$.
(iii) The proof of Lemma C.4.2 yields that for any $\mu \in \mathbb{C}$ and any $p \in \mathbb{N}$

$$
\mathfrak{A}^{j} \text { acts from } \mathscr{C}^{\infty}\left(\mathscr{E}, \mathfrak{H}_{\alpha, p}^{\mu}\right) \text { into } \mathscr{C}^{\infty}\left(\mathscr{E}, \mathfrak{H}_{\alpha, p}^{\mu+j-2 d}\right) \times \mathscr{C}^{\infty}\left(\mathscr{E}, \mathfrak{H}_{\alpha, p}^{\mu+j-\boldsymbol{\rho}}\right)
$$

Revisiting the proofs of Lemma C.4.1 and Proposition C.5.3 we obtain
Theorem C.7.2 Under the assumptions of Theorem C.6.2 and under hypothesis ( $\mathfrak{H}_{\mathrm{C} 1}$ ), for any $p \in \mathbb{N}_{0}$ the edge singular functions $\Phi^{p}$ belong to $\mathscr{C}^{\infty}\left(\mathscr{E}, \mathfrak{H}_{\alpha, p+1}^{\lambda+p}\right)$.

Since there holds for any $\mu \in \mathbb{C}$, for any $\alpha$ and $\alpha_{0} \in \mathbb{C}$

$$
\begin{aligned}
& (\alpha \zeta+\bar{\zeta})^{\mu}=\left(\alpha_{0} \zeta+\bar{\zeta}\right)^{\mu}+\sum_{k \geq 1} c_{\mu, k}\left(\alpha-\alpha_{0}\right)^{k} \zeta^{k}\left(\alpha_{0} \zeta+\bar{\zeta}\right)^{\mu-k} \\
& (\zeta+\alpha \bar{\zeta})^{\mu}=\left(\zeta+\alpha_{0} \bar{\zeta}\right)^{\mu}+\sum_{k \geq 1} c_{\mu, k}\left(\alpha-\alpha_{0}\right)^{k} \bar{\zeta}^{k}\left(\zeta+\alpha_{0} \bar{\zeta}\right)^{\mu-k}
\end{aligned}
$$

any function $\Phi$ in $\mathscr{C}^{\infty}\left(\mathscr{E}, \mathfrak{H}_{\alpha, p}^{\mu}\right)$ has a representation as

$$
\Phi=\sum_{\ell=1}^{m_{+}} \sum_{k=0}^{p n_{\ell}^{+}-1} \zeta^{k}\left(\alpha_{\ell}^{+}\left[x^{\prime}\right] \zeta+\bar{\zeta}\right)^{\mu-k} c_{k, \ell}^{+}\left[x^{\prime}\right]+\sum_{\ell=1}^{m_{-}} \sum_{k=0}^{p n_{\ell}^{-}-1} \bar{\zeta}^{k}\left(\zeta+\alpha_{\ell}^{-}\left[x^{\prime}\right] \bar{\zeta}\right)^{\mu-k} c_{k, \ell}^{-}\left[x^{\prime}\right],
$$

with $\mathscr{C}^{\infty}(\mathscr{E})$ coefficients $c_{k, \ell}^{ \pm}$. As a corollary of Theorem C.7.2 we obtain
Corollary C.7.3 Under the assumptions of Theorem C.6.2 and under hypothesis ( $\mathfrak{H}_{\mathrm{C} 1}$ ), for any $p \in \mathbb{N}_{0}$ the edge singular functions $\Phi^{p}$ have representations as

$$
\begin{align*}
\Phi^{p}\left[\mathscr{X}^{\prime}\right]= & \sum_{\ell=1}^{m_{+}} \sum_{k=0}^{(p+1) n_{\ell}^{+}-1} \zeta^{k}\left(\alpha_{\ell}^{+}\left[\mathscr{x}^{\prime}\right] \zeta+\bar{\zeta}\right)^{\lambda+p-k} \boldsymbol{c}_{k, \ell}^{p,+}\left[\mathscr{x}^{\prime}\right]  \tag{C.7.2}\\
& +\sum_{\ell=1}^{m_{-}} \sum_{k=0}^{(p+1) n_{\ell}^{-}-1} \bar{\zeta}^{k}\left(\zeta+\alpha_{\ell}^{-}\left[\mathscr{X}^{\prime}\right] \bar{\zeta}\right)^{\lambda+p-k} \boldsymbol{c}_{k, \ell}^{p,-}\left[\mathscr{x}^{\prime}\right], \quad \boldsymbol{c}_{k, \ell}^{p, \pm} \in \mathscr{C}^{\infty}(\mathscr{E})
\end{align*}
$$

Let us denote by $\Psi_{\ell, \omega}$ for $\omega= \pm 1$ the fundamental functions

$$
\Psi_{\ell,+}\left(x^{\prime}, r, \theta\right)=\alpha_{\ell}^{+}\left[x^{\prime}\right] \zeta+\bar{\zeta} \quad \text { and } \quad \Psi_{\ell,-}\left(x^{\prime}, r, \theta\right)=\zeta+\alpha_{\ell}^{-}\left[x^{\prime}\right] \bar{\zeta} .
$$

The comparison with the fundamental angular functions introduced in (B.8.3) is quite simple: Since, if $L$ is real,

$$
\tau_{\ell}=\frac{i\left(\alpha_{\ell}^{+}-1\right)}{\alpha_{\ell}^{+}+1} \quad \text { and } \quad \bar{\tau}_{\ell}=\frac{i\left(1-\alpha_{\ell}^{-}\right)}{\alpha_{\ell}^{-}+1}
$$

there holds

$$
\Psi_{\ell, \omega}(r, \theta)=\left(\alpha_{\ell}+1\right) r \psi_{\ell, \omega}(\theta), \quad \ell=1, \ldots, n_{\ell}^{\omega}, \quad \omega= \pm 1
$$

and conditions $\left(\mathfrak{H}_{\mathrm{B} 6}\right)$ and $\left(\mathfrak{H}_{\mathrm{C} 1}\right)$ are two formulations of the same assumption.
Coming back to the expansion (C.7.2), we note that the $N$-component vector functions

$$
\begin{aligned}
& \boldsymbol{d}_{\ell,+}^{p}\left(\mathscr{X}^{\prime}, y\right):=\sum_{k=0}^{(p+1) n_{\ell}^{+}-1} \zeta^{k} \Psi_{\ell,+}^{(p+1) n_{\ell}^{+}-1-k}\left(\mathscr{X}^{\prime}, y\right) \boldsymbol{c}_{k, \ell}^{p,+}\left[\mathscr{X}^{\prime}\right] \\
& \boldsymbol{d}_{\ell,-}^{p}\left(\mathscr{X}^{\prime}, y\right):=\sum_{k=0}^{(p+1) n_{\ell}^{-}-1} \bar{\zeta}^{k} \Psi_{\ell,-}^{(p+1) n_{\ell}^{-}-1-k}\left(\mathscr{X}^{\prime}, y\right) \boldsymbol{c}_{k, \ell}^{p,-}\left[\mathscr{X}^{\prime}\right]
\end{aligned}
$$

are polynomial in $y$, therefore $\mathscr{C}^{\infty}\left(\mathbb{R}^{n+1}\right)$, and there holds

$$
\begin{equation*}
\Phi^{p}\left[x^{\prime}\right]=\sum_{\omega= \pm 1} \sum_{\ell=1}^{m_{\omega}} \Psi_{\ell, \omega}^{\lambda-(p+1)\left(n_{\ell}^{\omega}-1\right)}\left(\mathscr{X}^{\prime}, y\right) \boldsymbol{d}_{\ell, \omega}^{p}\left(\mathscr{X}^{\prime}, y\right) \tag{C.7.3}
\end{equation*}
$$

If condition $\left(\mathfrak{H}_{\text {A3 }}\right)$ holds (i.e. if $\left.n_{\ell}^{ \pm}=1, \ell=1, \ldots, m\right)($ C.7.3) takes the simpler form

$$
\begin{equation*}
\Phi^{p}\left[\mathscr{X}^{\prime}\right]=\sum_{\omega= \pm 1} \sum_{\ell=1}^{m_{\omega}} \Psi_{\ell, \omega}^{\lambda}\left(\mathscr{X}^{\prime}, y\right) \boldsymbol{d}_{\ell, \omega}^{p}\left(\mathscr{X}^{\prime}, y\right) \tag{C.7.4}
\end{equation*}
$$

which means that the singular factors $\Psi_{\ell, \omega}^{\lambda}$ do not depend on $p$.
As a final consequence of formulas (C.7.3) and (C.7.4), we obtain "modular representations" of the solutions of elliptic BVP in the domain $\Omega=\mathbb{R}^{n+1} \backslash \mathscr{M}$ :
Theorem C.7.4 Let the hypotheses $\left(\mathfrak{H}_{\mathrm{A} 1}\right)$, $\left(\mathfrak{H}_{\mathrm{A} 2}\right)$ and $\left(\mathfrak{H}_{\mathrm{C} 1}\right)$ be satisfied.
(i) Any solution $\boldsymbol{u}$ of the boundary value problem (A.1.1) with smooth right hand side $\boldsymbol{f}$ has the following asymptotic expansion as $r \rightarrow 0$ : For any integer $K \geq 0$

$$
\begin{equation*}
\boldsymbol{u}=\sum_{\omega= \pm 1} \sum_{\ell=1}^{m_{\omega}} \Psi_{\ell, \omega}^{\frac{1}{2}-(K+1)\left(n_{\ell}^{\omega}-1\right)}\left(\mathscr{X}^{\prime}, y\right) \boldsymbol{d}_{\ell, \omega}^{[K]}\left(\mathscr{X}^{\prime}, y\right)+\boldsymbol{u}_{\mathrm{reg}, K}+\boldsymbol{u}_{\mathrm{rem}, K} \tag{C.7.5}
\end{equation*}
$$

where the vector-coefficients $\boldsymbol{d}_{\ell, \omega}^{[K]}$ are $\mathscr{C}^{\infty}\left(\mathbb{R}^{n+1}\right)$ and the regular parts $\boldsymbol{u}_{\mathrm{reg}, K}$ and $\boldsymbol{u}_{\mathrm{rem}, K}$ are as in Proposition A.4.2.
(ii) If the multiplicities $n_{\ell}^{\omega}$ are all equal to 1 , cf hypothesis $\left(\mathfrak{H}_{\mathrm{A} 3}\right)$, then $\boldsymbol{u}$ admits the global decomposition

$$
\begin{equation*}
\boldsymbol{u}=\sum_{\omega= \pm 1} \sum_{\ell=1}^{m} \Psi_{\ell, \omega}^{\frac{1}{2}}\left(x^{\prime}, y\right) \boldsymbol{d}_{\ell, \omega}^{\infty}\left(x^{\prime}, y\right)+\boldsymbol{u}_{\mathrm{reg}, \infty} \tag{C.7.6}
\end{equation*}
$$

where all vector-coefficients $\boldsymbol{d}_{\ell, \omega}^{\infty}$ and $\boldsymbol{u}_{\mathrm{reg}, \infty}$ are $\mathscr{C}^{\infty}\left(\mathbb{R}^{n+1}\right)$.

Remark C.7.5 The multiplicities $n_{\ell}^{\omega}$ are in fact the order of the poles of the inverse of the Cayley symbol $L^{\omega}(\alpha)^{-1}$ in $\alpha_{\ell}^{\omega}$. They can be smaller than the total multiplicity of $\alpha_{\ell}^{\omega}$. An example for this is the case of isotropic elasticity in $\mathbb{R}^{3}$ where $L^{ \pm}(\alpha)^{-1}$ have 0 as only pole, but the multiplicity is 2 (and not 3 ). The fundamental functions $\Psi_{\ell, \omega}$ are simply

$$
\Psi_{+}=\bar{\zeta}=\left(y_{1}-i y_{2}\right) \quad \text { and } \quad \Psi_{-}=\zeta=\left(y_{1}+i y_{2}\right)
$$

and expansion (C.7.5) takes the form, compare with [ChkDu2]

$$
\begin{equation*}
\boldsymbol{u}=\bar{\zeta}^{\frac{1}{2}-(K+1)} \boldsymbol{d}_{+}^{[K]}\left(x^{\prime}, y\right)+\zeta^{\frac{1}{2}-(K+1)} \boldsymbol{d}_{-}^{[K]}\left(\mathscr{x}^{\prime}, y\right)+\boldsymbol{u}_{\mathrm{reg}, K}+\boldsymbol{u}_{\mathrm{rem}, K} \tag{C.7.7}
\end{equation*}
$$

with $\mathscr{C}^{\infty}\left(\mathbb{R}^{n+1}\right)^{N}$ coefficients $\boldsymbol{d}_{ \pm}^{[K]}$.

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[^0]:    ${ }^{(1)}$ If the requirement $\boldsymbol{g} \in \mathscr{C}^{\infty}(\mathscr{M})^{N}$ is relaxed into $\boldsymbol{g} \in \mathrm{H}^{(\infty, s-\nu+K), \kappa}(\mathscr{M})^{N}$ for an integer $K>0$ and $\kappa \geq K$, we still obtain the asymptotics (B.4.3) for the same value $K$.

[^1]:    ${ }^{(2)}$ In [ChkDu1, Example 1.17] is considered the restriction of a $\Psi D O$ on $\mathbb{R}^{n+1}$ with a classical symbol onto the smooth surface $\mathscr{M}$ of codimension 1 and proved that the restricted operator is again a classical $\Psi D O$; explicit formulae for the components of the asymptotic expansion of the symbol are indicated.

[^2]:    ${ }^{(3)}$ For this, we use in particular that any contour integral of the integrand in (B.4.6) surrounding all the roots $\tau$ of $\operatorname{det} L\left(\mathscr{\mathscr { F }}_{\kappa}\left(\mathscr{E}^{\prime}\right)(\xi, \tau)\right)=0$, is zero, which allows to replace in (B.4.6) $\mathscr{L}_{+}$by $\mathscr{L}_{-}$.

[^3]:    ${ }^{(4)}$ Here we restrict consideration to the potential operators related to a second order system. For more general results we quote [ChkDu1] and the forthcoming paper: R.Duduchava, W.Wendland, Asymptotics of solutions to Agmon-Douglis-Nirenberg systems.

[^4]:    ${ }^{(5)}$ Defined by $\mathfrak{M}(\boldsymbol{f})\left[\mathscr{x}^{\prime}\right](\lambda, \theta)=\int_{0}^{\infty} r^{\lambda} \boldsymbol{f}\left(\mathscr{\mathscr { x }}^{\prime}, y\right) \frac{\mathrm{d} r}{r}$ as a natural extension of the formula on $\mathbb{R}_{+}$.
    ${ }^{(6)}$ More precisely, for any $K \in \mathbb{N}$ we obtain the contribution modulo $\mathscr{O}\left(r^{K}\right)$ to the infinite asymptotic series by using a contour which surrounds the (finite set of) poles $\lambda \in \Lambda\left[\mathscr{x}^{\prime}\right] \cup \mathbb{N}$ contained in the strip $d-1<\operatorname{Re} \lambda \leq K$.

[^5]:    ${ }^{(7)}$ The degree of homogeneity and the trace on $r=1$ completely determine an homogeneous function: If $v$ is homogeneous of degree $\mu$ and $v:=\left.v\right|_{r=1}$, then $v(r, \theta)=r^{\mu} V(\theta)$.

