# Richard B. Melrose <br> Antônio SÁ Barreto <br> Maciej Zworski <br> Semi-linear diffraction of conormal waves 

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## $\mathcal{N u m d a m}^{\prime}$

## 1996

## SEMI-LINEAR DIFFRACTION OF CONORMAL WAVES

Richard B. MELROSE, Antônio SÁ BARRETO and Maciej ZWORSKI

# SEMI-LINEAR DIFFRACTION OF CONORMAL WAVES 

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[^0]Résumé. - Nous étudions la régularité conormale de solutions bornées d'équations semi-linéaires strictement hyperboliques dans des domaines à bord diffractif:

$$
P u=f(x, u) \text { dans } X, u \upharpoonright_{\partial X}=0, u \in L_{\mathrm{loc}}^{\infty}(X)
$$

Si $X_{-} \subset X$ et $X$ est le domaine d'influence de $X_{-}$, nous considérons des solutions $u$ telles que $\operatorname{singsupp}(u) \cap X_{-} \cap \partial X=\emptyset$; de plus nous supposons que $u \upharpoonright_{X_{-}}$est conormale par rapport à une hypersurface caractéristique lisse, le front entrant.

Dans le cas de l'équation linéaire $f \equiv 0$, le support singulier de $u$ est contenu dans la réunion du front entrant et du front réfléchi obtenu par les lois de l'optique géométrique. Ces deux surfaces caractéristiques sont tangentes à l'ensemble des rayons rasants, le lieu des points où les bicaractéristiques entrantes sont tangentes au bord. Dans le cas semi-linéaire, nous démontrons que si de nouvelles singularités apparaissent alors elles apparaissent sur le demi-cône caractéristique au-dessus de l'ensemble des rayons rasants. En fait, le théorème de régularité conormale établi dans cet article est beaucoup plus précis.

Pour illustrer notre propos, nous choisirons pour $P$ l'opérateur des ondes à coefficients constants et pour $X$ le produit de $\mathbb{R}_{t}$ et de l'extérieur d'un obstacle strictement convexe. Alors $X_{-}=X \cap\{t<-T\}$. Comme donnée initiale, on pourra prendre une primitive locale de l'onde plane $\delta(t-\langle x, \omega\rangle)$ avec $T$ suffisamment grand. La géométrie de ce problème est figurée sur les schémas 1.1 et 1.2.

Abstract. - We study the conormal regularity of bounded solutions to semi-linear
strictly hyperbolic equations on domains with diffractive boundaries:

$$
P u=f(x, u) \text { in } X, u \upharpoonright_{\partial X}=0, u \in L_{\mathrm{loc}}^{\infty}(X) .
$$

If $X_{-} \subset X$ and $X$ is the domain of influence of $X_{-}$we consider solutions such that $\operatorname{singsupp}(u) \cap X_{-} \cap \partial X=\emptyset$ and further suppose that $u \upharpoonright_{X_{-}}$is conormal with respect to a smooth characteristic hypersurface, the incoming front.

For the linear equation, $f \equiv 0$, the singular support of $u$ is contained in the incoming front and the reflected front obtained using the rules of geometrical optics; these two characteristic surfaces are tangent at the glancing set, the locus of points at which the incoming bicharactersitics are tangent to the boundary. We prove that in the semi-linear case the only new singularites which may occur appear on the characteristic half-cone over the glancing set. The actual conormal regularity result presented in the paper is considerably more precise.

Our assumptions are best illustrated by taking for $P$ the constant coefficient wave equation with $X$ the product of $\mathbb{R}_{t}$ and the exterior of a strictly convex obstacle. Then $X_{-}=X \cap\{t<-T\}$ and for the initial data one can take locally an antiderivative of the plane wave $\delta(t-\langle x, \omega\rangle)$ with $T$ appropriately large. The geometry of this problem in two space dimensions is shown in Figures 1.1 and 1.2.

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## 1. INTRODUCTION AND STATEMENT OF RESULTS

The purpose of this paper is to describe the conormal regularity for a class of mixed problems for the semi-linear hyperbolic equations.

The study of $C^{\infty}$ regularity of solutions to non-linear wave equations has had two main directions: finding estimates on the strength of the anomalous singularities, i.e. those not present in the linear interaction, and obtaining geometric restrictions on the location of singularities. Our work is of the latter type. The strength of singularities for non-linear mixed problems has already been investigated with considerable success in $[45,10,21,48]$. The estimates on the location of singularities are much finer, so stronger assumptions are needed on the incoming waves or the initial data. The most striking example of this was provided by [2] where it is shown that wave-front set restrictions alone still allow the self-spreading of singularities, making the singular support propagate essentially in the same way as the support of the solution. Thus, in full generality, the location of singularities cannot be related to the original geometry except in a trivial way. A technically more challenging construction of a similar example for gliding mixed problems was then given in [47].

The appropriate class of distributions to consider for the incoming waves or the initial data are the conormal distributions, as was first noted in [6]. The conormal distributions appear naturally in the linear theory and are a subclass of the Lagrangian distributions motivated by geometrical optics. The interaction of conormal waves for interior problems has been investigated in [40, 32, 7, 9, 3, 42, 34] and the formation of non-linear caustics in $[18,19,11,27,43,44]$. For mixed problems, with only transversal reflections allowed, it was shown in [4, 5] that no anomalous singularities appear. One should also mention that examples of 'new' nonlinear singularities were provided at an early stage in [39]: namely, the interaction of three plane waves carrying conormal singularities produces a conic surface of new singularities propagating from the triple interaction point. However, in more complicated settings such as the propagation of the swallowtail or diffraction, where the 'new' cones are expected, no examples have yet been constructed. For the interior problems the methods developed in [20] provide a systematic approach to such constructions. Energy estimates used in the work on the lifespan of solutions to semi-linear hyperbolic equations $[15,22]$ are also, in essence, of conormal type.

If $\Sigma \subset X$ is a $C^{\infty}$ hypersurface in a $C^{\infty}$ manifold $X$, let ${ }^{\circ}(X, \Sigma)$ be the Lie algebra
of $C^{\infty}$ vector fields in $X$ tangent to $\Sigma$. The space of distributions of finite $L^{2}$-based conormal regularity with respect to $\Sigma$ is then defined by the stability of regularity under the applications of the elements of ()$(X, \Sigma)$ :

$$
I_{k} L_{\mathrm{loc}}^{2}(X, \Sigma)=\left\{u \in L_{\mathrm{loc}}^{2}(X): V_{1} \cdots V_{l} u \in L_{\mathrm{loc}}^{2}(X) \text { for } l \leq k \text { and } V_{i} \in \mathcal{O}(X, \Sigma)\right\}
$$

This modifies the definition of the Sobolev space $H_{(k)}$ by placing some geometric restrictions on the differentiations. Nevertheless, as observed in [32], bounded conormal functions have very good multiplicative properties in view of GagliardoNirenberg type inequalities.

Let us now consider a mixed hyperbolic problem with a diffractive boundary (see chapter 2 for a review of definitions). Our object of study is the semi-linear equation:

$$
\begin{equation*}
P u=f(x, u) \text { in } X,\left.u\right|_{\partial X}=0,\left.u\right|_{X_{-}}=u_{0} \tag{1.1}
\end{equation*}
$$

where $f$ is a $C^{\infty}$ function of its arguments, $P$ is a strictly hyperbolic operator, $X$ is a $C^{\infty}$ manifold with the boundary $\partial X, X_{-}=\{x \in X: \phi(x)<-T\}$ with $\phi \in C^{\infty}(X)$ a time function for $P$ and the time $T$ fixed.

The initial data is assumed to be conormal to the incident front $F$. The reflection rule of geometrical optics produces the reflected front $R$. With the motivation coming again from the geometrical optics we define the shadow boundary on $\partial X$ as

$$
\Gamma=\partial X \cap \operatorname{cl}[R \cap F \backslash \partial X]
$$

The front obtained from the nonlinear interaction is the forward half-cone, $S_{+}$, of $P$-bicharacteristics starting on $\Gamma$. Let us also denote by $D_{+}$and $B_{+}$the two components of the set of glancing characteristics on $S_{+}$. A more detailed discussion of the fronts is presented in chapter 2. Fig. 1.1 shows three different time slices and Fig. 1.2 is a space-time picture. Note that $R$ and $F$ are hypersurfaces with singular boundaries.

The crudest form of our result is
Theorem 1.1. - Let $u \in L^{\infty}(X)$ be a bounded solution of (1.1) with

$$
u_{0} \in I_{\infty} L_{\mathrm{loc}}^{2}\left(X_{-} ; F\right)
$$

Then

$$
W F_{b}(u) \subset{ }^{b} N^{*} R \cup{ }^{b} N^{*} F \cup{ }^{b} N^{*} S_{+} \cup{ }^{b} N^{*} B_{+} \cup{ }^{b} N^{*} D_{+} \cup{ }^{b} T_{\Gamma}^{*} X \backslash 0
$$

We refer the reader to [25] and [14], Sect. 18.3 for the definition of the $b$-wave front set, $W F_{b}$, which reduces to the ordinary $W F$ away from the boundary $\partial X$. We use the natural map $\jmath: T^{*} X \backslash 0 \rightarrow{ }^{b} T^{*} X \backslash 0$ (see chapter 4 and the references given above) to define ${ }^{b} N^{*} \Sigma=\jmath\left(N^{*} \Sigma\right)$.

Theorem 1.1 immediately gives the singular support statement:
Corollary 1.2. - Under the assumptions of Theorem 1.1

$$
\text { sing supp } u \subset F \cup R \cup S_{+} .
$$



Figure 1.1. The fronts projected to the space variables at fixed times

Since the data $u_{0}$ is conormal, one would like to describe precisely the conormal regularity of the solution $u$. In fact the proof is based on the construction of an appropriate space with good multiplicative and propagative properties - see chapter 3. Since the precise definition of this 'strong', but not quite conormal, space is rather involved we shall content ourselves with a weaker statement here, referring the reader to Definition 3.2 and Theorem 8.2 for the full result.

Theorem 1.3. - Let $u \in L_{\mathrm{loc}}^{\infty}(X)$ be a bounded solution of (1.1) with

$$
u_{0} \in I_{k} L_{\mathrm{loc}}^{2}\left(X_{-} ; F\right)
$$

If $\Omega$ is an open subset of $X$ such that

$$
\Omega \cap\left(D_{+} \cup B_{+}\right)=\emptyset
$$

then

$$
\left.u\right|_{\Omega} \in I_{k} L_{\mathrm{loc}}^{2}(\Omega, F)+I_{k} L_{\mathrm{loc}}^{2}(\Omega, R)+I_{k} L_{\mathrm{loc}}^{2}\left(\Omega, S_{+}\right) .
$$

Already in the transversal case this is slightly stronger than the result in [4] as conormal singularities with respect to the boundary are excluded.

Our conclusions are concerned purely with the $L^{2}$-based regularity. The present existence theory [45] requires higher Sobolev regularity for $u_{0}$ to guarantee local existence of bounded solutions, so one needs to assume $u_{0} \in I_{k} L_{\mathrm{loc}}^{2}\left(X_{-} ; F\right) \cap H_{(s)}\left(X_{-}\right)$ for $s>n / 2$. However, the conormal results described above should lead to an improvement in the style of [41]. It should be noted that our present method does not treat the fully semi-linear equation $P u=f(x, u, \nabla u)$, essentially because the iteration procedure in $k$ proceeds in steps of $1 / 2$.


Figure 1.2. The forward half-cone and the glancing boundaries $B$ and $D$.

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## 2. DIFFRACTIVE GEOMETRY

First we describe the interaction of a characteristic hypersurface for a second-order hyperbolic operator with a bicharacteristically concave (diffractive) boundary. In particular the reflected front is shown to have a cusp singularity when continued across the boundary.

Let $X$ be a manifold with boundary equipped with a pseudo-Riemannian metric of hyperbolic signature,,,,$+---\ldots$ The metric symbol $p \in S^{2}\left(T^{*} X\right)$ is therefore a polynomial of degree two on each fibre and it can be reduced, in linear coordinates in each fibre, to

$$
\tau^{2}-\xi_{1}^{2}-\cdots-\xi_{n}^{2}, \quad \operatorname{dim} X=n+1
$$

The boundary of $X$ is said to be time-like if $p$ is negative-definite on $N^{*} \partial X$; this is always assumed below. It will be convenient to assume that $X$ is time oriented; this amounts to the continuous selection of one of the solid cones, $p>0$, in the fibres. A function $t \in C^{\infty}(X)$ is a time-function if $p(d t)>0$.

The assumption that $\partial X$ is time-like means that it carries an induced pseudoRiemannian metric of hyperbolic signature. If $g$ is the dual quadratic form to $p$, on $T X$, then $g_{\partial}=g \upharpoonright_{T \partial X}$ fixes the induced structure. Let $p_{\partial}$ denote the metric symbol on $T^{*} \partial X$. Set

$$
\left.\begin{array}{l}
\mathscr{H}=\left\{p_{\partial}>0\right\}=\mathscr{H}_{+} \cup \mathscr{H}_{-} \\
\mathscr{G}=\left\{p_{\partial}=0\right\} \\
\mathscr{G}=\left\{p_{\partial}<0\right\},
\end{array}\right\} \quad \text { in } T^{*} \partial X \backslash 0
$$

respectively the hyperbolic, glancing and elliptic regions of $T^{*} \partial X \backslash 0$. The timeorientation of $X$ induces a time-orientation of $\partial X$, giving the decomposition of the hyperbolic region.

The restriction to the boundary of the characteristic variety

$$
\Sigma=\{p=0\} \subset T^{*} X \backslash 0
$$

projects onto $\mathscr{H} \cup \mathscr{G}=\overline{\mathscr{H}}:$

$$
\begin{equation*}
\Sigma_{\partial}=\Sigma \cap T_{\partial X}^{*} X \xrightarrow{\iota^{*}} \mathscr{H} \cup \mathscr{G} . \tag{2.1}
\end{equation*}
$$

Here $\iota^{*}: T_{\partial X}^{*} X \longrightarrow T^{*} \partial X$ is the pull-back map induced by the inclusion $\partial X \stackrel{\iota}{\hookrightarrow} X$.

The map (2.1) is a fold, $2-1$ over $\mathscr{H}$ and $1-1$ over $\mathcal{G}$. To see this, note that $\iota^{*}: T_{\partial X}^{*} X \longrightarrow T^{*} \partial X$ is the quotient map with respect to $N^{*} \partial X$. The restriction of $p$ to each fibre of $\iota^{*}$ is of this form

$$
\begin{equation*}
p=\iota^{*} p_{\partial}-n^{2} \tag{2.2}
\end{equation*}
$$

where $n: T_{\partial X}^{*} X \longrightarrow \mathbb{R}$ is the symbol of the inward-pointing vector field, $V$, satisfying $g(V, \cdot)=0$ on $T \partial X, g(V, V)=-1$. Thus $p_{\partial} \geq 0$ on the range of $\iota^{*}$ on $\Sigma_{\partial}$ and the involution exchanging the points in $\Sigma_{\partial}$ with the same image satisfies

$$
\begin{equation*}
I_{\partial}: \Sigma_{\partial} \longleftrightarrow \Sigma_{\partial}, I_{\partial}^{*} n=-n, \iota^{*} \cdot I_{\partial}=\iota^{*} \tag{2.3}
\end{equation*}
$$

Let $\tilde{\mathscr{H}}_{ \pm}, \tilde{\mathcal{E}}$ be the preimages of $\mathscr{H}_{ \pm}$and $\mathscr{G}$ under $\iota^{*}$, so in particular $\iota^{*}: \tilde{\mathcal{E}} \longleftrightarrow \mathscr{S}$ is an isomorphism.

The projection $\iota^{*}$ can be expressed symplectically. Let $x \in C^{\infty}(X)$ be a defining function for $\partial X$ and let $q \in C^{\infty}\left(T^{*} X\right)$ be its lift to $T^{*} X$. Then the Hamilton vector field $H_{q}$ satisfies $H_{q} q=0$, i.e. $H_{q}$ is tangent to the leaves of $\iota^{*}: T_{\partial X}^{*} X \longrightarrow T^{*} \partial X$. Since the leaves are tangent to $\Sigma$ exactly at $\tilde{\mathscr{S}}$,

$$
\tilde{\mathscr{G}}=\{p=0, q=0,\{p, q\}=0\}
$$

where $\{p, q\}=H_{p} q=-H_{q} p$ is the Poisson bracket. The simple tangency of $H_{q}$ to $\Sigma_{\partial}$, corresponding to the fact that (2.1) is a fold, is expressed by

$$
\{q,\{q, p\}\}<0
$$

This holds throughout $T_{\partial X}^{*} X$, since $\partial X$ is time-like. Applying $H_{q}$ to both sides of (2.2) and noting that $\{q, n\}$ is the lift of a function from the base, so $\{q,\{q, n\}\} \equiv 0$

$$
\begin{gathered}
\{q, p\}=-2 n\{q, n\} \\
\{q,\{q, p\}\}=-2\{q, n\}^{2}
\end{gathered}
$$

Thus (2.2) can be written in terms of Poisson brackets

$$
\begin{equation*}
p=\iota^{*} p_{\partial}+\frac{\{q, p\}^{2}}{2\{q,\{q, p\}\}} \tag{2.4}
\end{equation*}
$$

The denominator $\{q,\{q, p\}\}$ is also the lift of a function from the base. Thus, if the involution is extended to $T_{\partial X}^{*} X$ so that the second two conditions in (2.3) hold, then

$$
\begin{equation*}
I_{\partial}^{*}\{q, p\}=-\{q, p\} \tag{2.5}
\end{equation*}
$$

The points of $\mathscr{G}$ are further distinguished by the behaviour of the second Poisson bracket:

$$
\begin{align*}
& \tilde{\mathscr{G}}_{d}=\left\{m \in \tilde{\mathscr{G}} ; H_{p}^{2} q(m)>0\right\} \\
& \tilde{\mathscr{G}}_{h}=\left\{m \in \tilde{\mathscr{G}} ; H_{p}^{2} q(m)=0\right\}  \tag{2.6}\\
& \tilde{\mathfrak{G}}_{g}=\left\{m \in \tilde{\mathscr{S}} ; H_{p}^{2} q(m)<0\right\}
\end{align*}
$$

These, and similarly their images in $\mathcal{G}$ under $\iota^{*}$, are respectively the sets of diffractive, higher-order and gliding points. The boundary of $X$ is said to be diffractive (or bicharacteristically concave) if $\mathscr{G}=\mathcal{G}_{d}$; this is always assumed below.

Consider the differential of the involution $I_{\partial}$ on $\Sigma_{\partial}$, at $\mathcal{G}$. In particular note that the type of a glancing point, in the sense of (2.6), is reflected in:

$$
\begin{equation*}
\left(I_{\partial}\right)_{*} H_{p}-H_{p}=2 \frac{\{p,\{p, q\}\}}{\{q,\{q, p\}\}} H_{q} \tag{2.7}
\end{equation*}
$$

In fact, since $\iota^{*} \cdot I_{\partial}=I_{\partial}$, the projection of $\left(I_{\partial}\right)_{*} H_{p}$ under $\iota^{*}$ is equal to the projection of $H_{p}$. The null space of the projection is spanned by $H_{q}$, so it is only necessary to compute the coefficient on the right. Applying both sides to $\{q, p\}$ and using (2.5) gives (2.7).

We will be concerned with the local geometry near a base point $x_{0} \in \partial X$, so we are free to shrink $X$ as necessary. In this sense the assumption that the boundary is diffractive is really that $\mathcal{G} \cap T_{x_{0}}^{*} \partial X \subset \mathcal{G}_{d}$. In case $X=\mathbb{R} \times Y$ carries a product metric, $g=d t^{2}-h$, with $h$ a Riemannian metric on $Y$, the boundary is diffractive if and only if $\partial Y$ is strictly geodesically concave. In case $Y=\mathbb{R}^{n} \backslash K$ where $K$ is an open, smoothly bounded region and $h$ is the Euclidean metric this is equivalent to the strict convexity of $K$ (cf. [26]).

It is convenient to consider an extension, $\tilde{X}$, of $X$ to a manifold without boundary. A corresponding extension of this pseudo-Riemannian structure will be denoted $\tilde{p}$. The defining function $x \in C^{\infty}(X)$ extends to $\tilde{x} \in C^{\infty}(\tilde{X})$ and if $\widetilde{X}$ is chosen small enough, $\partial X=\{\tilde{x}=0\}$ is an embedded hypersurface. The freedom to shrink $X$ will be used to choose $\widetilde{X}$ to be bicharacteristically convex.

In $\widetilde{X}$ we consider a closed characteristic hypersurface for $\tilde{p}$, passing through this point $x_{0}$. Thus $\widetilde{F} \subset \widetilde{X}$ satisfies

$$
\begin{equation*}
\tilde{F}=\{\tilde{f}=0\}, \quad \tilde{f} \in C^{\infty}(\tilde{X}), d \tilde{f} \neq 0 \text { on } \tilde{F},\left.\tilde{p}\right|_{N^{*} \tilde{F}}=0 \tag{2.8}
\end{equation*}
$$

Since $\partial X$ is time-like, $N^{*} \partial X$ and $N^{*} \widetilde{F}$ are linearly independent and hence

$$
\widetilde{F}_{\partial}=\tilde{F} \cup \partial X \hookrightarrow \partial X
$$

is an embedded hypersurface. Since $N^{*} \widetilde{F}_{\partial}=\iota^{*}\left(N_{\partial X}^{*} \widetilde{F}\right)$, with $N_{\partial X}^{*} \widetilde{F} \subset \Sigma_{\partial}$, we have $N_{\widetilde{F}}^{*} \widetilde{F}_{\partial} \subset \overline{\mathcal{H}}=\mathscr{H} \cup \mathcal{G}$. For us the most interesting points are the diffractive points for $\widetilde{F}$ :

Lemma 2.1. - If $\widetilde{F} \subset \widetilde{X}$ is an embedded characteristic hypersurface then

$$
\begin{equation*}
\Gamma=\pi\left(N^{*} \widetilde{F}_{\partial} \cap \mathscr{G}_{d}\right) \subset \widetilde{F}_{\partial} \tag{2.9}
\end{equation*}
$$

is an embedded hypersurface.
Proof. - Consider the function $\tilde{\gamma}=\{p, q\}$ restricted to $N^{*} \tilde{F}$. By assumption $N^{*} \widetilde{F} \subset \Sigma$ so $H_{p}$ is tangent to it. At any point of $\tilde{\mathscr{G}}_{d} \cap N_{\partial X}^{*} \widetilde{F},\{p, q\}=0$ and
$\{p,\{p, q\}\}>0$, by definition. Thus $H_{p}\{p, q\} \neq 0$ and so $d \tilde{\gamma} \neq 0$ on $N^{*} \widetilde{F}$. Since $H_{p} q=\{p, q\}=0$ on $\tilde{\mathscr{G}}_{d}, q$ and $\{p, q\}$ have independent differentials on $N^{*} \widetilde{F}$ at $\tilde{\mathcal{G}}_{d} \cap N_{\partial X}^{*} \widetilde{F}$. It follows that

$$
N_{\partial X}^{*} \widetilde{F} \cap \mathcal{G}_{d} \subset N_{\partial X}^{*} \widetilde{F}
$$

is an embedded hypersurface. Since it is homogeneous it projects to the embedded hypersurface $\Gamma$.

The characteristic hypersurface $\widetilde{F}$ is to be thought of as the extension through the boundary of $X$ of the incident front. It is important to separate which parts of $\widetilde{F}$ are intrinsic and which depend on the choice of extension-the latter being necessarily irrelevant to the final form of the results.

By assumption $N^{*} \widetilde{F}$ is closed, so

$$
\Lambda_{F} \stackrel{\text { def }}{=} N^{*} \widetilde{F} \backslash 0
$$

is the union of the maximally extended bicharacteristic intervals, i.e. integral curves of $H_{p}$, through each of its points. Set

$$
F=\left\{z \in \widetilde{F} \cap X ; \text { the bicharacteristics through } N_{z}^{*} \widetilde{F} \text { stay in } T^{*} X \text { for } t \leq t(z)\right\}
$$

Here, $t$ is a time function. The submanifold $\Gamma \subset F$ is the singular locus in $F$ near which it is not even a manifold with corners. Indeed the boundary of $F$ consists of two smooth manifolds with boundary (each of codimension two in $X$ )

$$
\begin{equation*}
\partial F=F_{\partial} \cup B, F_{\partial} \cap B=\partial F_{\partial}=\partial B=B \cap \partial X=\Gamma . \tag{2.10}
\end{equation*}
$$

Here $F_{\partial}$ is half of $\widetilde{F}_{\partial}$ and $B$, the shadow boundary, is the projection into $X$ of the forward half-bicharacteristic starting at points of $N_{\Gamma}^{*} \widetilde{F}$.

The main objective of this section is to consider the reflected front generated by $\widetilde{F}$ and $\partial X$. To do so we need to recall the notion of a hypersurface with cusp singularity. By definition a cusp hypersurface is one which is diffeomorphic to $C=\left\{x_{2}^{3}=x_{1}^{2}\right\}$ in $\mathbb{R}^{n}, n \geq 2$.

A simple characterization can be obtained in terms of the closure of the conormal bundle to the regular part of the hypersurface. As is easily checked

$$
\Lambda_{C}=\operatorname{cl} N^{*}\left\{x_{2}^{3}=x_{1}^{2} ; x_{2}>0\right\} \subset T^{*} \mathbb{R}^{n} \backslash 0
$$

is a smooth, homogeneous Lagrangian. Now a point of the singular locus, $L=\left\{x_{1}=\right.$ $\left.x_{2}=0\right\}$,

$$
\begin{align*}
& \pi: \Lambda_{C} \longrightarrow \mathbb{R}^{n} \text { has differential with } \\
& \text { two-dimensional null space at } \Lambda_{C} \cap T_{l}^{*} \mathbb{R}^{n}, l \in L . \tag{2.11}
\end{align*}
$$

Moreover, any vector field $V$ on $T^{*} \mathbb{R}^{n}$ which is tangent to $T_{l}^{*} \mathbb{R}^{n}$, independent of the radial vector field and takes the value $v \in T_{m} \Lambda_{C} \cap T_{m}\left(T_{l}^{*} \mathbb{R}^{n}\right)$ at $m$ is only simply
tangent to $\Lambda_{C}$ at $m$. Conversely (see Arnol'd [1]) if these two conditions hold for $\Lambda_{C}$ near $m \in T_{l}^{*} \mathbb{R}^{n} \cap \Lambda_{C}$ then the projection of a neighbourhood of $m \in \Lambda_{C}$ is a cusp. We use this abstract characterization, with $\mathbb{R}^{n}$ replaced by $\widetilde{X}$ (as can obviously be done) to analyze the reflected front.

Set $\widetilde{\Lambda}_{R}^{0}=I_{\partial}\left(N_{\partial X}^{*} \widetilde{F}\right)$ and let $\widetilde{\Lambda}_{R}$ be the $H_{p}$-flow-out in $T^{*} \widetilde{X} \backslash 0$ of $\widetilde{\Lambda}_{R}^{0}$. Thus $\widetilde{\Lambda}_{R}$ is just the union of the maximally extended $H_{p}$ integral curves passing through points of $\widetilde{\Lambda}_{R}^{0}$.

Proposition 2.2. - If $\tilde{F} \subset \tilde{X}$ is a smooth characteristic hypersurface for which $x_{0} \in$ $\partial X$ is a diffractive point then, for $\widetilde{X}$ shrunk to a sufficiently small bicharacteristically convex neighbourhood of $x_{0}, \widetilde{\Lambda}_{R} \subset T^{*} \widetilde{X} \backslash 0$ is a smooth closed conic Lagrangian submanifold which is the closure of the conormal bundle to a hypersurface with cusp singularity, $\widetilde{R}$, through $x_{0}$.

Proof. - To see that $\widetilde{\Lambda}_{R}$ is smooth it suffices to observe that $H_{p}$ is not tangent to the initial surface $\widetilde{\Lambda}_{R}^{0}$. By definition $\widetilde{\Lambda}_{R}^{0}$ is the image under $I_{\partial}$ of $N_{\partial X}^{*} \widetilde{F}$ and $N_{x_{0}}^{*} \widetilde{F} \backslash 0 \subset \mathcal{G}_{d}$, by assumption. Thus $H_{p}$ is tangent to $N_{\partial X}^{*} \widetilde{F}$ at $N_{x_{0}}^{*} \widetilde{F} \backslash 0$ and hence $\left(I_{\partial}\right)_{*} H_{p}$ is tangent to $\widetilde{\Lambda}_{R}^{0}$ at $N_{x_{0}}^{*} \widetilde{F}$. Now $H_{q}$ cannot be tangent to $N_{\partial X}^{*} \widetilde{F} \backslash 0$ (since this would mean $\widetilde{F}$ was tangent to $\partial X$ ) and hence cannot be tangent to $\widetilde{\Lambda}_{R}^{0}$ at $N_{x_{0}}^{*} \widetilde{F}$. From (2.7) it follows that $H_{p}$ is not tangent to $\widetilde{\Lambda}_{R}^{0}$, so $\widetilde{\Lambda}_{R} \subset T^{*} \widetilde{X} \backslash 0$ is smooth, closed and conic if $\tilde{X}$ is chosen small enough. It is also invariant under reflection in the fibres.
This discussion shows that both $\left(I_{\partial}\right)_{*} H_{p}$ and $H_{p}$ are tangent to $\widetilde{\Lambda}_{R}$ at $N_{x_{0}}^{*} \widetilde{F}$, hence so is $H_{q}$. Since $H_{q}$ is tangent to the fibres of $T^{*} \widetilde{X}$ and is non-radial at $N_{x_{0}}^{*} \widetilde{F} \backslash 0$ it follows that the differential of the projection $\pi: \widetilde{\Lambda}_{R} \longrightarrow \widetilde{X}$ has null space of dimension of at least two. In fact $\pi_{*}: T_{m} \widetilde{\Lambda}_{R} \longrightarrow T_{\pi(m)} \widetilde{X}$ has rank exactly $\operatorname{dim} \widetilde{X}-2$, since $\pi_{*}: T_{m} \widetilde{\Lambda}_{R}^{0} \longrightarrow T_{\pi(m)} \partial X$ has rank $\operatorname{dim} \partial X-1$ at $m \in N_{x_{0}}^{*} \widetilde{F} \backslash 0$. Finally note that $H_{q}$ is only simply tangent to $\widetilde{\Lambda}_{R}$ at $N_{x_{0}}^{*} \widetilde{F} \backslash 0$ since it is only simply tangent to $\Sigma$. Thus $\widetilde{\Lambda}_{R}$ is the closure of the conormal bundle of a hypersurface with cusp singularity

Clearly the cusp locus $L \subset \widetilde{R}$ passes through $\Gamma$. It is important to check that

$$
\begin{equation*}
L \backslash \Gamma \subset \tilde{X} \backslash X \text { and } L \text { is simply tangent to } \partial X \text { at } \Gamma . \tag{2.12}
\end{equation*}
$$

Since the tangent space to $L$ is just the image of the tangent space to $\widetilde{\Lambda}_{R}$ under the projection, $L$ is certainly tangent to $\partial X$ at $\Gamma$. The reflected front $\widetilde{R} \cap X^{\circ}$ is smooth so the inclusion follows. To see the simple tangency we first choose coordinates so that $\tilde{X} \subset \mathbb{R}_{x}^{3} \times \mathbb{R}_{y}^{n-3}$ and

$$
\widetilde{R}=\left\{(x, y): x_{3}^{2}=x_{2}^{3}\right\} \cap \widetilde{X}, \quad \Gamma=\{(x, y): y=0\} \cap \widetilde{X}
$$

As in Proposition 2.1 of [34] it then follows that

$$
\begin{equation*}
p=a\left(4 \xi_{2}^{2}-9 x_{2} \xi_{3}^{2}+\xi_{1} p_{1}+\langle\eta, \tilde{p}\rangle\right), a \neq 0 \tag{2.13}
\end{equation*}
$$

The cusp locus is given by $L=\left\{x_{2}=x_{3}=0\right\} \cap \tilde{X}$ and if $\partial X=\{\rho=0\}, X=\{\rho>0\}$ then the tangency of $L$ to $\partial X$ at $\Gamma$ implies

$$
\begin{equation*}
\partial_{x_{1}} \rho(0, y)=0 \tag{2.14}
\end{equation*}
$$

The points $m=(0, y ;(0,0,1), 0) \in T^{*} \tilde{X} \backslash 0$ are diffractive:

$$
p(m)=\rho(m)=\{p, \rho\}(m)=0, \quad\{p,\{p, \rho\}\}(m)>0
$$

Hence (2.13) and (2.14) show that

$$
\begin{equation*}
\partial_{x_{2}} \rho(0, y)>0 \tag{2.15}
\end{equation*}
$$

and consequently

$$
\rho(x, y)=x_{2}-g\left(x_{1}, x_{3}, y\right), \quad \partial_{x_{1}} g(0, y)=g(0, y)=0 .
$$

Since $L \backslash \Gamma \subset\{\rho<0\}$ we also see that $g\left(x_{1}, 0, y\right) \geq 0$ and thus we can write

$$
\begin{equation*}
g\left(x_{1}, x_{3}, y\right)=g_{1}\left(x_{1}, y\right) x_{1}^{2}+x_{3} g_{2}\left(x_{1}, x_{3}, y\right) \tag{2.16}
\end{equation*}
$$

The restriction of $\widetilde{R}$ to $\partial X$ is given by

$$
\operatorname{cl}(\partial X \cap \widetilde{R} \backslash \widetilde{F}) \cup(\partial X \cap \widetilde{F})=\widetilde{R} \cap \partial X=\left\{(x, y) \in \partial X: g\left(x_{1}, x_{3}, y\right)^{3}=x_{3}^{2}\right\}
$$

and using (2.16) we see that

$$
\begin{aligned}
\widetilde{F} \cap \partial X= & \left\{(x, y): x_{3}=G\left(x_{1}, x_{3}, y\right) g_{1}\left(x_{1}, y\right)^{3}+F\left(x_{1}, x_{3}, y\right) g_{1}\left(x_{1}, y\right)^{4},\right. \\
& \left.x_{2}=g\left(x_{1}, x_{3}, y\right)\right\}, G(0, y) \neq 0, \\
\operatorname{cl}(\widetilde{R} \backslash \widetilde{F}) \cap \partial= & \left\{(x, y): x_{3}=-G\left(x_{1}, x_{3}, y\right) g_{1}\left(x_{1}, y\right)^{3}+F\left(x_{1}, x_{3}, y\right) g_{1}\left(x_{1}, y\right)^{4}\right. \\
& \left.x_{2}=g\left(x_{1}, x_{3}, y\right)\right\} .
\end{aligned}
$$

The Lagrangian $N^{*}(\widetilde{F} \cap \partial X)$ is simply tangent to $\mathcal{G} \subset T^{*} \partial X \backslash 0$ and since $N^{*}(\operatorname{cl}(\widetilde{R} \backslash \widetilde{F}) \cap \partial X)$ is related to it by the billiard ball map (see (2.19)) below the two Lagrangians are simply tangent at $N^{*}(\widetilde{F} \cap \partial X) \cap \mathcal{G}$ (see the equivalent model case in chapter 7 ). Hence $\partial_{x_{1}} g_{1}(0, y) \neq 0$ and thus in view of (2.16), $L$ is simply tangent to $\partial X=\left\{x_{2}=g\left(x_{1}, x_{3}, y\right)\right\}$.

In the case of the wave equation in the exterior of a convex obstacle Proposition 2.1 was given in [49]. In that case the cusp locus $L$ projected to the space variables is the envelope of the reflected rays, see Fig. 2.1.

In the remainder of this section we review the geometry of the fronts. With the possible exception of Proposition 2.4 b ) (see [34]) all the facts are essentially well known and are implicit in the proofs in chapter 3 and chapter 7 below.

The intersection properties of $\widetilde{F}$ and $\widetilde{R}$ are described in
Proposition 2.3. - Let $B$ be the shadow boundary defined by (2.10). Then

## SEMI-LINEAR DIFFRACTION OF CONORMAL WAVES



Figure 2.1. The extended reflected front projected to the space variables at a fixed time.

1. $\widetilde{F} \cap \widetilde{R}=(\widetilde{F} \cap \partial X) \cup B$
2. $\Lambda_{R}$ and $\Lambda_{F}$ intersect cleanly and $\Lambda_{R} \cap \Lambda_{F}=T_{B}^{*} X \backslash 0 \cap \Lambda_{R}=T_{B}^{*} X \backslash 0 \cap \Lambda_{F}$.

It is important to remark that although the extension $\tilde{p}$ was used in the definition of $\Lambda_{R}$, the part of $\widetilde{R}$ corresponding to the true reflection is determined by $p$ and $F$ alone. It will be denoted by $R$ and is defined as follows. Proposition 2.3 implies that $\widetilde{R} \backslash \widetilde{F}$ has four components, two of which are disjoint from $L$. We now take as $R$ the closure of the one for which $R \cap \partial X=F \cap \partial X$. An alternative definition is provided by taking
$R=\left\{z \in \widetilde{R} \cap X: \iota^{*}\left(B_{z}^{-} \cap T_{\partial X}^{*} X \backslash 0\right) \subset \iota^{*}\left(N_{\partial X}^{*} F\right)\right.$, where $B_{z}^{-}$is the bicharacteristic through $N_{z}^{*} \widetilde{R}$ with $\left.t \leq t(z)\right\}$.

The bicharacteristic cone over the shadow boundary in $\partial X, \Gamma$ is now defined in the standard way, as the union of the maximally extended bicharacteristic intervals over $N^{*} \Gamma \cap \Sigma$. We denote it by $\Lambda_{S}$ and its projection by $\widetilde{S}$. We note however that $\left.\widetilde{S}\right|_{X}$ depends on the extension $\tilde{p}$. Thus we need

Proposition 2.4. - The set $D=\operatorname{cl}\left[\pi_{*}\left(\exp H_{p}\left(N^{*} \Gamma \cap \mathcal{G}\right)\right) \backslash B\right]$ is a smooth codimension two submanifold of $X$, tangent to $\partial X$ at $\Gamma$ and intersecting $B$ cleanly, $B \cap D=\Gamma$. The components $S_{+}$and $S_{-}$of $\widetilde{S} \backslash(B \cup D)$ such that $S_{ \pm} \cap(\tilde{X} \backslash X)=\emptyset$ are determined by $\Gamma$ and $p$.

We can separate the forward and retarded components, $S_{ \pm}$, respectively, and similarly denote by $\widetilde{S}_{+}$the full forward cone over $\Gamma$. We also denote by $B_{+}$and $D_{+}$the intersections of $B$ and $D$ with $\mathrm{cl} S_{+}$respectively, see Fig. 3.1

In the non-linear interaction more geometry is present. In addition to the cone over $\Gamma$ we will also have to include, in a very residual way, smooth characteristic surfaces tangent to $\widetilde{S}$ at $D$ - see chapter 3 . Thus we define

$$
\begin{equation*}
\mathscr{R}=\left\{H \subset \widetilde{X} \text { smooth hypersurface : }\left.p\right|_{N^{*} H}=0, N_{D \backslash \Gamma}^{*} H=N_{D \backslash \Gamma}^{*} \widetilde{S}\right\} \tag{2.17}
\end{equation*}
$$

with the first easy observation that $\bigcap_{H \in \mathcal{K}} H=D$. By the analogy with the previous notation we also write $\Lambda_{H} \stackrel{\text { def }}{=} N^{*} H \backslash 0$.

The intersection properties of $\Lambda_{S}$ and $\Lambda_{R}, \Lambda_{F}, \Lambda_{H}$ are given in

## Proposition 2.5. - The following Lagrangian submanifolds

1. $\Lambda_{S}$ and $\Lambda_{F}$ intersect cleanly at $\Lambda_{R} \cap \Lambda_{S}=\Lambda_{S} \cap N^{*} B$,
2. $\Lambda_{S}$ and $\Lambda_{R}$ are simply tangent along $\Lambda_{R} \cap \Lambda_{F}=\Lambda_{S} \cap N^{*} B$,
3. $\Lambda_{S}$ and $\Lambda_{H}$ intersect cleanly at $\Lambda_{S} \cap N^{*} D$.

The restrictions of the characteristic (singular) surfaces $\widetilde{R}$ and $\widetilde{S}$ to the boundary can be related to $F_{\partial}=\widetilde{F} \cap \partial X$ and $\Gamma \subset \partial X$ respectively. Thus we define the following smooth Lagrangian submanifolds of $T^{*} \partial X \backslash 0$ :

$$
\begin{equation*}
\Lambda_{00}=N^{*} \Gamma, \Lambda_{11}=N^{*}(\widetilde{F} \cap \partial X), \Lambda_{31}=N^{*}(\operatorname{cl}(\widetilde{R} \backslash \tilde{F}) \cap \partial X), \Lambda_{13}=N^{*}(H \cap \partial X) \tag{2.18}
\end{equation*}
$$

and

$$
\Lambda_{21} \cup \Lambda_{23}=N^{*}(\widetilde{S} \cap \partial X), \quad \Lambda_{11} \cap \Lambda_{23}=\emptyset, \quad \Lambda_{13} \cap \Lambda_{21}=\emptyset
$$

Using the time function restricted to the boundary, $t \upharpoonright \partial X$, we then obtain Lagrangians with boundaries:

$$
\Lambda_{i j}^{ \pm}=\Lambda_{i j} \cap\{ \pm t \geq 0\}, \quad j=1,3, i=1,2,3 .
$$

The relations between these Lagrangians in then given by the billiard ball map $\delta^{ \pm}: \imath^{*}\left(\Sigma_{\partial}\right) \rightarrow \imath^{*}\left(\Sigma_{\partial}\right)$ (see $\left.[35,36]\right)$ which is two valued $( \pm)$ and has a square root singularity over $\mathscr{G}$. To recall the definition of $\delta$ we introduce the natural map

$$
\gamma_{p}: \Sigma_{\partial} \longrightarrow \Sigma / H_{p}
$$

which, same as $\imath^{*}$ above, has a simple fold. Thus we can associate to it an involution $I_{p}$ exchanging the points with the same image (as $I_{\partial}$ did for $\imath^{*}$ ). The two valued map $\delta^{ \pm}$is then defined by

$$
\delta^{ \pm} \circ \imath^{*}=I_{p}
$$

The relation between the Lagrangians with boundaries is then given by

$$
\begin{equation*}
\Lambda_{31}^{ \pm}=\delta^{ \pm} \Lambda_{11}^{ \pm}, \quad \Lambda_{21}^{ \pm} \cup \Lambda_{23}^{ \pm}=\delta^{ \pm} \Lambda_{0} \tag{2.19}
\end{equation*}
$$

Finally, we relate the geometry described above to the differential equation. We assume that the principal symbol of a strictly hyperbolic operator $P$ in $X, p$, satisfies the glancing assumptions with respect to the boundary $\partial X$ defined by $q$. On the extension $\tilde{X}$, we introduce a strictly hyperbolic extension of $P$ with the symbol $\tilde{p}$. We will also denote it by $P$ keeping in mind, however, the freedom we have in its choice. The past $\widetilde{X}_{-} \subset \widetilde{X}$ is defined so that $\widetilde{X}$ is its domain of influence, and

$$
\tilde{F} \cap \tilde{X}_{-}=F \cap \tilde{X}_{-} \subset X \cap \tilde{X}_{-} \stackrel{\text { def }}{=} X_{-}, \quad F \cap \tilde{X}_{-} \cap \partial X=\emptyset
$$

## 3. RESOLUTION OF SINGULARITIES AND THE CONORMAL SPACES

As described in chapter 2 the interaction geometry is quite complicated as it involves cusp and conic singularities. To define a conormal space with reasonable propagation of regularity for $P$, one follows the method originating from [32] and subsequently applied in $[27,34,42,43,44]$. Its essence is the resolution of singularities and the use of the vector fields tangent to the lifted geometry in the resolved space. The insistence on conormality is motivated by the good multiplicative properties of bounded conormal functions, as already indicated in chapter 1 and the conviction that conormal regularity excludes any hidden singularities that could produce selfspreading.

For our problem the method of resolution is similar to that used in [34, 44] and it involves a non-homogeneous blow-up. To describe it let us consider

$$
\mathbb{R}^{n}=\mathbb{R}^{3} \times \mathbb{R}^{n-3}, \quad z=(x, y), z \in \mathbb{R}^{n}, x \in \mathbb{R}^{3}, y \in \mathbb{R}^{n-3}
$$

on which we define an $\mathbb{R}_{+}-$action $T_{\delta}^{1-2-3}$ :

$$
\begin{equation*}
T_{\delta}^{1-2-3}(x, y)=\left(\delta x_{1}, \delta^{2} x_{2}, \delta^{3} x_{3}, y\right), \quad \delta \in \mathbb{R}_{+} \tag{3.1}
\end{equation*}
$$

We start with a definition of spaces of functions with given non-homogeneous orders of vanishing:

$$
\begin{equation*}
M_{r}^{1-2-3}(\tilde{X}) \subset C^{\infty}(\tilde{X}), \quad u \in M_{r}^{1-2-3}(\tilde{X}) \Longleftrightarrow T_{\delta}^{*} u=\mathcal{O}\left(\delta^{r}\right), \delta \rightarrow 0 \tag{3.2}
\end{equation*}
$$

This allows us to the define a filtration of the differential operators in terms of homogeneity. Thus

$$
Q \in \operatorname{Diff}_{p,(1-2-3)}^{k}(\tilde{X}) \Longleftrightarrow Q: M_{r}^{1-2-3}(\tilde{X}) \longrightarrow M_{r-p}^{1-2-3}(\widetilde{X}) \quad \text { for } \quad r \geq p
$$

The homogeneous differential operator important in our discussion is Friedlander's operator in $\mathbb{R}^{3}$ :

$$
P_{0}=4 D_{x_{2}}^{2}-9 x_{2} D_{x_{3}}^{2}-6 D_{x_{3}} D_{x_{1}} .
$$

Proposition 3.1. - There exist coordinates $(x, y) \in \mathbb{R}^{n}, x \in \mathbb{R}^{3}, y \in \mathbb{R}^{n-3}$ in $\tilde{X}$ such that

$$
\begin{gather*}
P=P_{0}+Q, \quad Q \in \operatorname{Diff}_{3,(1-2-3)}^{2}(\tilde{X})  \tag{3.3}\\
\partial X=\left\{(x, y): x_{2}-\frac{1}{4} x_{1}^{2}=0\right\} \tag{3.4}
\end{gather*}
$$

and with the notation of chapter 2 and any $H \in \mathcal{R}$ given by (2.17)

1. $\Gamma=\left\{(0, y): y \in \mathbb{R}^{n-3}\right\} \cap \tilde{X}$
2. $\widetilde{R}=\left\{(x, y): x_{3}^{2}-x_{2}^{3}=0, x \in \mathbb{R}^{3}, y \in \mathbb{R}^{n-3}\right\}$
3. $\widetilde{F}=\left\{(x, y): 2 x_{3}-3 x_{2} x_{1}+x_{1}^{3}+f_{4}=0, x \in \mathbb{R}^{3}, y \in \mathbb{R}^{n-3}\right\}, f_{4} \in M_{4}^{1-2-3}\left(\mathbb{R}^{n}\right)$
4. $\widetilde{S}=\left\{(x, y): x_{1}^{4}+8 x_{1} x_{3}-6 x_{1}^{2} x_{2}-3 x_{2}^{2}+s_{5}=0, x \in \mathbb{R}^{3}, y \in \mathbb{R}^{n-3}\right\}$, $s_{5} \in M_{5}^{1-2-3}\left(\mathbb{R}^{n}\right)$.
5. $H=\left\{(x, y): x_{1}+h_{2}=0\right\}, \quad h_{2} \in M_{2}^{1-2-3}\left(\mathbb{R}^{n}\right)$.

Proof. - By Proposition 2.1, $\widetilde{R}$ is a surface with a cusp singularity and thus we can find coordinates $(x, y)$ such that $\widetilde{R}$ is given by $b$ ). The cusp locus is then $L=\left\{x_{2}=x_{3}=0\right\}$ and as noted in (2.12) it is simply tangent to $\partial X$ at $\Gamma$, $L \backslash \Gamma \subset \widetilde{X} \backslash X$. If $X$ is given by $\rho(x, y)>0$ near $(0,0)$ then by $(2.15) \rho_{x_{2}}^{\prime}(0,0)>0$. We conclude that

$$
\rho(x, y)=x_{2}-g\left(x_{1}, x_{3}, y\right), \quad \partial_{x_{1}} g(0, y)=0
$$

near $(0,0)$. We expand $g$ into $g_{0}\left(x_{3}, y\right)+x_{1} g_{1}\left(x_{3}, y\right)+x_{1}^{2} g_{2}\left(x_{1}, x_{3}, y\right)$ where $g_{0}(0, y)=$ $g_{1}(0, y)=0$ and $g_{2}>0$. (since $L \backslash \Gamma \subset \widetilde{X} \backslash X$ ). Completing the square we write $g$ as

$$
\left(g_{0}\left(x_{3}, y\right)-\frac{1}{16} \frac{\left(g_{1}\left(x_{3}, y\right)\right)^{2}}{g_{2}\left(x_{1}, x_{3}, y\right)}\right)+\frac{1}{4}\left(2 x_{1}\left(g_{2}\left(x_{1}, x_{3}, y\right)\right)^{\frac{1}{2}}+\frac{1}{4} \frac{g_{1}\left(x_{3}, y\right)}{\left(g_{2}\left(x_{1}, x_{3}, y\right)\right)^{\frac{1}{2}}}\right)^{2}
$$

We then change variables by replacing $x_{1}$ by the term in the second bracket and observe that the term in the first bracket, say $\tilde{g}$, vanishes for all $x_{1}$ when $x_{3}=0$. We can now apply Theorem 4.1 of [1] to obtain a cusp preserving change of variables $\left(x_{2}, x_{3}\right)$ depending smoothly on the parameters $\left(x_{1}, y\right)$ and putting $x_{2}-\tilde{g}\left(x_{1}, x_{3}, y\right)$ to $x_{2}$. Thus in the new coordinates $\rho(x, y)=x_{2}-x_{1}^{2} / 4$ as desired.
If $\widetilde{R}=\left\{x_{3}^{2}-x_{2}^{3}=0\right\}$ is characteristic for a strictly hyperbolic operator $P$ then one easily sees that $P=P_{0}+Q, Q \in \mathrm{Diff}_{3,1-2-3}^{2}$ - see Proposition 2.1, [34]. We observe also that neglecting $Q$ produces errors of higher homogeneity (see the proof of Proposition 2.1 [34].). Since the right hand side of $c$ ) is the reflection of $\widetilde{R}$ with respect to $\partial X$ obtained using the symbol of $P_{0}$, we conclude that
$\widetilde{F}=\{(x, y): f(x, y)=0\}, \quad f(x, y)=2 x_{3}-3 x_{2} x_{1}+x_{1}^{3}+f_{4}\left(x_{1}, x_{2}, y\right), f_{4} \in M_{4}^{1-2-3}$.

The same conclusion can be made about $\widetilde{S}$ :

$$
\begin{equation*}
\widetilde{S}=\{(x, y): q(x, y)=0\}, \quad q(x, y)=x_{1}^{4}+8 x_{1} x_{3}-6 x_{1}^{2} x_{2}-3 x_{2}^{2}+q_{5}, q_{5} \in M_{5}^{1-2-3} \tag{3.6}
\end{equation*}
$$

The surface $B$ is given by $\widetilde{S} \cap \widetilde{F} \cap \widetilde{R}$ which up to terms of higher homogeneity is parametrized by $\left(t, t^{2}, t^{3}, y\right)$. The codimension two surface $D$ consists of characteristics of $P$ tangent to $\partial X$ and not contained in $B$. Neglecting $Q$ in $P, D$ one would obtain $D=\{(0,0, t, y)\}$ (which would lie in $\partial X$ ). Since that describes $D$ up to terms of higher homogeneity we conclude that $D=\left\{\left(f_{1}(t, y), f_{2}(t, y), t, y\right)|t|<\epsilon\right\}$ and it easily follows from (3.6) that $\nabla q \upharpoonright_{D}=t\left[\left(8, \tilde{f}_{1}(t, y), \widetilde{f}_{2}(t, y)+O(t)\right], \widetilde{f_{i}} \in C^{\infty}\right.$. Since $H$ is smooth and tangent to $\widetilde{S}$ at $D$, it follows that its normal at $D$ is given by $\left(8, \widetilde{f}_{1}, \widetilde{f}_{2}\right)+O(t)$. The implicit function theorem immediately gives $e$ ) with $h=h\left(x_{2}, x_{3}, y\right)$ satisfying $h(0,0, y)=0$, that is $h \in M_{2}^{1-2-3}\left(\mathbb{R}^{n}\right)$.

We will consider the surfaces on the right-hand side as the model geometry. The sense in which they are models can be explained as follows. The model surface for $\widetilde{F}$ in $c$ ) is characteristic for Friedlander's operator $P_{0}$ and the cusp $\widetilde{R}$ is obtained from that model surface by reflection (according to the rules of geometric optics given in Proposition 2.1) through the boundary $x_{2}-\frac{1}{4} x_{1}^{2}=0$. Note that this surface, although microlocally diffractive near $N^{*} \widetilde{R}$, is not globally diffractive for $P_{0}$ : it contains the characteristic $\left\{x_{1}=x_{2}=0\right\}$. Thus we see that $Q \neq 0$ and essentially it has to contain a term of the form $-c x_{2} D_{x_{1}}^{2}$ which destroys the degeneracy of the characteristic $\left\{x_{1}=x_{2}=0\right\}$. The surface defined by the right hand side of $d$ ) is the cone over $0 \in \mathbb{R}^{3}$ with respect to the characteristic flow-out by $P_{0}$.

In view of Proposition 3.1 it is natural to resolve the geometry using the 1-2-3 blow-up given by the $\mathbb{R}_{+}$-action (3.1). Thus we define the space

$$
\begin{equation*}
\tilde{X}_{1}=(\tilde{X} \backslash \Gamma) \sqcup\left(\mathbb{S}_{1-2-3}^{2} \times \mathbb{R}^{n-3}\right) \simeq \mathbb{R}_{+} \times \mathbb{S}_{1-2-3}^{2} \times \mathbb{R}^{n-3} \tag{3.7}
\end{equation*}
$$

where $\mathbb{S}_{1-2-3}^{2}$ is a non-round sphere $\left\{\omega \in \mathbb{R}^{3}: \sum_{1 \leq i \leq 3} \omega_{i}^{12 / i}=1\right\}$ and where the $C^{\infty}$ structure on $\widetilde{X}_{1}$ is given by the second identification (see [28]). We now have the blow-down map

$$
\widetilde{X}_{1} \xrightarrow{\beta_{1}} \widetilde{X}, \quad(r, \omega, y) \longmapsto\left(r \omega_{1}, r^{2} \omega_{2}, r^{3} \omega_{3}, y\right)
$$

which is a diffeomorphism on $\widetilde{X}_{1} \backslash \partial \widetilde{X}_{1}$. Thus following [28] we define the pull-back of $Y$ to be

$$
\beta_{1}^{*} Y=\operatorname{cl}\left[\beta_{1}^{-1}(Y \backslash \Gamma)\right], \quad Y \subset \tilde{X}
$$

Propositions 3.1 imply that $\beta_{1}^{*} \widetilde{F}, \beta_{1}^{*} \widetilde{S}$ and $\beta_{1}^{*} \partial X$ are smooth hypersurfaces in $\tilde{X}_{1}$ intersecting the boundary $\partial \widetilde{X}_{1}$ cleanly, and $\beta_{1}^{*} \widetilde{R}$ has a cusp singularity transversal to $\partial \widetilde{X}_{1}$. Also,

$$
\begin{aligned}
& \beta_{1}^{*} \widetilde{F} \cap \partial \widetilde{X}_{1}=\beta_{1}^{*}\left\{2 x_{3}-3 x_{2} x_{1}+x_{1}^{3}=0\right\} \cap \partial \widetilde{X}_{1}, \\
& \beta_{1}^{*} \widetilde{S} \cap \partial \widetilde{X}_{1}=\beta_{1}^{*}\left\{x_{1}^{4}+8 x_{1} x_{3}-6 x_{1}^{2} x_{2}-3 x_{2}^{2}=0\right\} \cap \partial \widetilde{X}_{1}, \\
& \beta_{1}^{*} H \cap \partial \widetilde{X}_{1}=\beta_{1}^{*}\left\{x_{1}=0\right\} \cap \partial \widetilde{X}_{1} .
\end{aligned}
$$

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The boundary of the resolved space $\partial \widetilde{X}_{1}$ and the above intersections are shown in Fig. 3.1.


Figure 3.1. The geometry on $\partial \widetilde{X}_{1}$ as seen from the positive $x_{2}$ direction.
The lift of the operator $P$ is of the form:

$$
r^{4} \beta_{1_{*}} P \beta_{1}^{*}=P_{1} \in \operatorname{Diff}_{b}^{2}\left(\widetilde{X}_{1}\right), \quad{ }^{b} \sigma_{2}\left(P_{1}\right) \upharpoonright^{b} T_{\partial \widetilde{X}_{1}}^{*} \widetilde{X}_{1}={ }^{b} \sigma_{2}\left(r^{4} \beta_{1_{*}} P_{0} \beta_{1}^{*}\right) \upharpoonright^{b} T_{\partial \widetilde{X}_{1}}^{*} \widetilde{X}_{1}
$$

where we refer the reader to [25] (and also [14], Sect. 18.3) for the definition of the totally characteristic operators, the compressed cotangent bundle and the $b$-symbol map.

As the blow-ups described above depend on the particular choice of coordinates, or alternatively the $\mathbb{R}_{+}$-actions, it is important to note certain invariance properties. To state them, let $\chi: \widetilde{X} \rightarrow \widetilde{X}$ be a local diffeomorphism such that $\chi(\Gamma) \subset \Gamma$. For simplicity we shall consider the 1-2-3 action given by $T_{\delta}=T_{\delta}^{1-2-3}$ only and require that

$$
\begin{equation*}
\left[I d-\left(T_{\delta} \circ \chi\right)^{-1} \circ\left(\chi \circ T_{\delta}\right)\right]^{*}: M_{r}^{1-2-3} \longrightarrow M_{r+1}^{1-2-3}, \quad \text { for all } r \in \mathbb{N}_{0} \tag{3.8}
\end{equation*}
$$

This almost homogeneity condition is now present in
Proposition 3.2. - Any diffeomorphism $\chi$ satisfying (3.8) lifts to a boundary preserving diffeomorphism

$$
\chi_{1}: \tilde{X}_{1} \longrightarrow \widetilde{X}_{1}
$$

Also, if a diffeomorphism preserves $\widetilde{R}$ and $\widetilde{S}$ in coordinates given by Proposition 3.1, then it satisfies (3.8).

For the proof and further geometric results needed in the study of interior propagation we refer the reader to [34]. In chapter 6 we will however need the following result:

Proposition 3.3. - If $\rho$ is a defining function of $\partial X$ then each of the following can be obtained with a diffeomorphism satisfying (3.8):
i) $H=\left\{x_{1}=0\right\}, \widetilde{S}=\left\{x_{2}^{2}-4 x_{1} x_{3}=0\right\}, B=\left\{x_{3}=x_{2}=0\right\}, d \rho \upharpoonright_{\Gamma}=d x_{2}$,
ii) $\widetilde{F}=\left\{2 x_{3}-3 x_{1} x_{2}+x_{1}^{3}=0\right\}, \widetilde{R}=\left\{x_{3}^{2}-x_{2}^{3}=0\right\}, d \rho \upharpoonright_{\Gamma}=d x_{2}+\alpha d x_{3}$.

Proof. - A diffeomorphism leading to $i$ ) is already essentially obtained in Proposition 3.3 of [34] with the almost homogeneity (3.8) guaranteed by construction. This implies that

$$
d \rho \upharpoonright_{\Gamma}=\alpha d x_{2}+\beta d x_{3}, \quad \alpha \neq 0
$$

If $\beta \neq 0$, then $D=\widetilde{S} \cap H=\left\{x_{2}=x_{1}=0\right\}$ would be transversal to the boundary which contradicts the glancing assumption. Thus, $\beta=0$ and by rescaling $\alpha=1$.
For $i i$ ) we first recall from Proposition $3.1 b$ ) that $\widetilde{F}$ is given by $2 x_{3}-3 x_{1} x_{2}+x_{1}^{3}+$ $h_{4}\left(x_{1}, x_{2}, y\right)=0, h_{4} \in M_{4}^{1-2-3}$. Also, $B=\operatorname{cl}(\widetilde{R} \cap \widetilde{F} \backslash X)$ is up to terms of higher homogeneity given by the normal form $\left(t, t^{2}, t^{3}, y\right)$, that is, for small $t$ :

$$
B=\left\{\left(t, t^{2}(1+t \phi(t, y))^{2}, t^{3}(1+t \phi(t, y))^{3}, y\right)\right\}
$$

where we used the exact form of $\widetilde{R}$. Thus changing $x_{1}$ to $x_{1}\left(1+x_{1} \phi\left(x_{1}, y\right)\right)$ puts $B$ into its normal form. Since the defining function of $\underset{\widetilde{F}}{\widetilde{F}}$ has to vanish there and its gradient has to coincide with the normal of $\widetilde{R}(\widetilde{R}$ and $\widetilde{F}$ are simply tangent along $B$ ) we conclude that

$$
h_{4}\left(t, t^{2}, y\right)=0, \partial_{x_{1}} h_{4}\left(t, t^{2}, y\right)=\partial_{x_{2}} h_{4}\left(t, t^{2}, y\right)=0
$$

and thus $h_{4}\left(x_{1}, x_{2}, y\right)=\left(x_{2}-x_{1}^{2}\right)^{2} \tilde{h}_{4}\left(x_{1}, x_{2}, y\right)$. A simple applications of the homotopy method concludes the proof. In fact, let us put $f_{s}=2 x_{3}-3 x_{2} x_{1}+x_{1}^{3}+s h_{4}$. We want to find a cusp ( $\widetilde{R}$ ) preserving family of diffeomorphisms such that

$$
\chi_{s}^{*} f_{s}=f_{0}, 0 \leq s \leq 1
$$

If $V_{s}$ is the family of vector fields generating $\chi_{s}$, the required conditions are equivalent to

$$
V_{s} h_{s}=-h_{4}, \quad V_{s} \text { tangent to } \widetilde{R}
$$

This is obtained near $\Gamma$ with

$$
V_{s}=\frac{\left(x_{2}-x_{1}^{2}\right) \tilde{h}_{4}}{3+4 s z \tilde{h}_{4}+s\left(x_{2}-x_{1}^{2}\right)\left(\tilde{h}_{4}\right)_{x_{1}}^{\prime}} \frac{\partial}{\partial x_{1}} .
$$

Since the diffeomorphism constructed above clearly satisfies (3.8) we conclude that $\rho \equiv x_{2}-x_{1}^{2} / 4\left(\bmod M_{3}^{1-2-3}\right)$.

Remark 3.4. - In case i) considered in the proposition

$$
\rho(x, y)=x_{2}+c_{1}(y) x_{1}^{2}+c_{2}(y) x_{2}^{2}+c_{3}(y) x_{3}^{2}+\sum_{i<j} a_{i j}(y) x_{i} x_{j}+\mathcal{O}\left(|x|^{3}\right)
$$

and the assumption that $\left\{x_{1}=x_{2}=0\right\}$ and $\left\{x_{2}=x_{3}=0\right\}$ are glancing easily implies that $c_{1}, c_{3}>0$.

The map $\Psi\left(x_{1}, x_{2}, x_{3}\right)=\left(2 x_{3}+x_{1}^{3}-3 x_{2} x_{1}, x_{2}-x_{1}^{2}, \frac{4}{3} x_{1}, y\right)$ transforms $Q$ to $\left\{4 x_{3} x_{1}-x_{2}^{2}=0\right\}$, thus the cone on the right hand side of $d$ ) in Proposition 3.1 is essentially symmetric with respect to the interchange of $x_{1}$ and $x_{3}$. Roughly speaking, an additional blow-up near $\beta_{1}^{*}(Q \cap G) \cap \partial \widetilde{X}_{1}$ is needed to undo the asymmetry of the 1-2-3 blow-up.

To introduce it we first change coordinates (using a diffeomorphism satisfying (3.8)) so that $i$ ) of Proposition 3.3 holds. Using the lift of these coordinates, we blow-up with the 2-1-1 homogeneity the codimension three submanifold $\partial \widetilde{X}_{1} \cap \beta_{1}^{*}\left\{x_{1}=x_{2}=\right.$ $\left.0, x_{3}>0\right\}=\partial \widetilde{X}_{1} \cap \beta_{1}^{*} D_{+}:$

$$
\begin{gathered}
\tilde{X}_{2} \xrightarrow{\beta_{12}} \tilde{X}_{1} \xrightarrow{\beta_{1}} \tilde{X}, \quad \beta_{2}=\beta_{1} \circ \beta_{12} \\
\widetilde{X}_{2}=\widetilde{X}_{1} \backslash\left(\partial \widetilde{X}_{1} \cap \beta_{1}^{*} D_{+}\right) \sqcup\left(\mathbb{S}_{2-1-1_{+}}^{2} \times \mathbb{R}^{n-3}\right),
\end{gathered}
$$

where $\mathbb{S}_{2-1-1_{+}}^{2}$ is a half non-round sphere $\left\{\nu \in \mathbb{R}^{3}: \nu_{1}^{2}+\nu_{2}^{4}+\nu_{3}^{4}=1, \nu_{3} \geq 0\right\}$ and

$$
\beta_{12}(\rho, \nu, y)=\left(\rho^{2} \nu_{1}, \rho \nu_{2}, \rho \nu_{3}, y\right)
$$

with the coordinates in $\widetilde{X}_{1}$ near $\partial \widetilde{X}_{1} \cap \beta_{1}^{*} D_{+}$chosen so that

$$
\beta_{1}\left(X_{1}, X_{2}, r, y\right)=\left(r X_{1}, r^{2} X_{2}, r^{3}, y\right) \in \tilde{X}
$$

The manifold $\widetilde{X}_{2}$ has a codimension two corner and $\partial \widetilde{X}_{2}$ is shown in Fig. 3.2.


Figure 3.2. The hierarchy of blow-ups

Since $\widetilde{S}$ and $H$ are simply tangent at $D$ another blow-up is still needed:

$$
\tilde{X}_{3} \xrightarrow{\beta_{23}} \tilde{X}_{2} \xrightarrow{\beta_{12}} \widetilde{X}_{1} \xrightarrow{\beta_{1}} \tilde{X}, \quad \beta_{3}=\beta_{1} \circ \beta_{12} \circ \beta_{23} .
$$

Here, the line $\beta_{2}^{*} D_{+}$is blown-up with the 2-1-0 homogeneity in the coordinates where $\beta_{2}^{*} H \cap N=\left\{X_{1}=0\right\} \cap N$ and $\beta_{2}^{*} \widetilde{S} \cap N=\left\{4 X_{1}-X_{2}^{2}=0\right\} \cap N$ where $N$ is a neighbourhood of $\beta_{2}^{*} D_{+}$- see [43].

There are additional tangencies and singularities that have not yet been resolved: the tangencies described in Proposition 2.5 persist in $\widetilde{X}_{1}$ at $\beta_{1}^{*} B$ as does the cusp singularity of $\beta_{1}^{*} \widetilde{R}$ at $\beta_{1}^{*} L$. The former is resolved using a succession of normal blowups [33] (see Fig. 3.3) and the latter using the 3-2 blow up [43], only at $\beta_{1}^{*} B_{+}$and $\beta_{1}^{*} L_{+}$respectively $\left(\beta_{1}^{*} L_{+}=\beta_{1}^{*}\left(L \cap\left\{x_{1} \geq 0\right\}\right.\right.$. This leads to the space $\widetilde{X}_{4}$ :

$$
\tilde{X}_{4} \xrightarrow{\beta_{34}} \widetilde{X}_{3} \xrightarrow{\beta_{3}} \tilde{X}, \quad \beta=\beta_{4}=\beta_{3} \circ \beta_{34}
$$

see Fig. 3.2.
For future reference we also define $\widetilde{X}_{5}$, analogously to $\widetilde{X}_{4}$ but obtained by applying the same blow-ups at the lifts of $D_{ \pm}, L_{ \pm}, B_{ \pm}$rather than of $D_{+}, L_{+}, B_{+}$only:

$$
\tilde{X}_{5} \xrightarrow{\beta_{35}} \tilde{X}_{3} \xrightarrow{\beta_{3}} \tilde{X}, \quad \beta_{5}=\beta_{3} \circ \beta_{35} .
$$

We shall now define the $C^{\infty}$-algebra $J_{k} L_{c}^{2}(\tilde{X}, H)$ associated to the geometry in the open manifold $\tilde{X}$. In the notation we stress the dependence on the 'artificial' characteristic hypersurface $H \in \mathscr{R}$.

Let us first consider the surfaces in $\tilde{X}_{4}$ obtained from the geometry in $\tilde{X}$ :

$$
\beta^{*} \widetilde{F}, \beta^{*} \widetilde{R}, \beta^{*} \widetilde{S}_{+}, \beta^{*}(\widetilde{F} \cap \widetilde{R} \backslash B), \beta^{*}\left(\widetilde{S}_{+} \cap \widetilde{R} \backslash B\right), \beta^{*} H
$$

where we note that the lifts of $B_{+}, D_{+}$and $L_{+}$are included in the boundary of $\widetilde{X}_{4}$. Let $\mathcal{S}$ be the variety obtained by taking a disjoint union of the five submanifolds above with $\partial \widetilde{X}_{4}$ :

$$
\begin{equation*}
\mathcal{S}=\beta^{*} \widetilde{F} \sqcup \beta^{*} \widetilde{R} \sqcup \beta^{*} \widetilde{S}_{+} \sqcup \beta^{*}(\widetilde{F} \cap \widetilde{R} \backslash B) \sqcup \beta^{*}\left(\widetilde{S}_{+} \cap \widetilde{R} \backslash B\right) \sqcup \beta^{*} H \sqcup \partial \widetilde{X}_{4} \tag{3.9}
\end{equation*}
$$

Ideally, we would want to define $J_{k} L_{c}^{2}(\widetilde{X}, H)$ as the $\beta$-pushforward of the conormal spaces associated to $\mathscr{S}$ which is in fact done for the interior problem. Here, however, this would be disastrous.

In chapter 5 we shall define $K_{1}=K_{1}(\epsilon) \subset \widetilde{X}_{1}$ (see Fig. 5.1) which in some sense constitutes a 'non-homogeneous' past. We can take $\epsilon$ small enough so that

$$
\begin{equation*}
\beta_{1}^{*}(\widetilde{F} \cap \widetilde{R} \backslash B) \subset \widetilde{X}_{1} \backslash K_{1}(\epsilon) \Leftrightarrow \beta_{1}^{*}(\widetilde{F} \cap \partial X) \subset \widetilde{X}_{1} \backslash K_{1}(\epsilon) \tag{3.10}
\end{equation*}
$$

Since the all the higher generation blow-ups occur away from $K_{1}$ we can think of it as a subset of $\widetilde{X}_{4}\left(\right.$ or $\left.\beta_{14}^{*} K_{1}=K_{1}\right)$.


Figure 3.3. The three normal blow-ups of $\beta_{1}^{*} B$

Definition 3.5. - For $k \in \mathbb{N}_{0}$, we define

$$
J_{k} L_{c}^{2}(\tilde{X}, H)=\beta_{*}\left\{U \in I_{k} L_{\nu}^{2}\left(\tilde{X}_{4} ; \mathcal{O}\left(\tilde{X}_{4}, \mathscr{\delta}\right)\right): U \upharpoonright_{K_{1}} \in I_{k} L_{\nu}^{2}\left(\tilde{X}_{4}, \beta^{*} \tilde{F} \sqcup \partial \widetilde{X}_{4}\right) \upharpoonright_{K_{1}}\right\}
$$

where the variety $\mathfrak{5}$ is given by (3.9) and $K_{1}=K_{1}(\epsilon)$ is given by Definition 5.1 with $\epsilon$ such that (3.10) is satisfied. The norm is defined using the norm of the lift:

$$
\|u\|_{J_{k} L^{2}(\widetilde{X}, H)}=\left\|\beta^{*} u\right\|_{I_{k} L_{\nu}^{2}\left(\widetilde{X}_{4}, \tilde{\nu}\left(\widetilde{X}_{4}, \tilde{\mathfrak{s}}\right)\right)}+\left\|\left.\beta^{*} u\right|_{K_{1}}\right\|_{I_{k} L_{\nu}^{2}\left(\widetilde{X}_{4}, \beta^{*} \widetilde{F} \sqcup \partial \widetilde{X}_{4}\right) \Gamma_{K_{1}}}
$$

We also define $J_{k}^{1} L^{2}(\tilde{X}, H)$ by demanding that $u, D u \in J_{k} L_{c}^{2}(\tilde{X}, H)$, with the obvious norm.

For non-integral values of the order of regularity we use complex interpolation and define:

$$
\begin{equation*}
J_{k+s} L^{2}(\tilde{X}, H)=\left[J_{k} L^{2}(\tilde{X}, H), J_{k+1} L^{2}(\tilde{X}, H)\right]_{s}, \quad 0<s<1 \tag{3.11}
\end{equation*}
$$

and similarly for $J_{k+s}^{1} L^{2}(\widetilde{X}, H)$.
This is a pseudo-conormal space as it involves an additional condition in $K_{1}$. The corresponding pseudo-conormal space for the manifold with boundary $X$ is essentially obtained by restriction with an additional singular support condition:

Definition 3.6. - For $s \geq 0, s \in \mathbb{R}$ we define

$$
\begin{aligned}
J_{s} L^{2}(X)= & \left\{u \in L^{2}(X): \text { for every } H \in \mathscr{R} \text { there exists } \tilde{u} \in J_{s} L^{2}(\widetilde{X}, H)\right. \\
& \text { with } \tilde{u}\left\lceil_{X}=u \text { and sing supp }{ }^{(s)} u \cap\left(\widetilde{F} \backslash F \cup \widetilde{R} \backslash R \cup \widetilde{S}_{+} \backslash S_{+}\right)=\emptyset\right\}
\end{aligned}
$$

We recall that using the regularity function $s_{u}(x)$ (cf. [14], Sect. 18.1), we define sing supp ${ }^{(s)} u=\left\{x: s_{u}(x) \leq s\right\}$ which by lower semi-continuity of $s_{u}$ is closed. The space $J_{s} L^{2}(X)$ is not a normed space and although it can be made into a Fréchet space we shall not need this fact here. The $L_{\text {loc }}^{2}$ based spaces are defined in the obvious way: $u \in J_{s} L_{\text {loc }}^{2}(X)$ if and only if for any $\chi \in C_{0}^{\infty}(X), \chi u \in J_{s} L^{2}(X)$.

Remark 3.7. - Although the definition of the blow-up involves the choice of $H$, it can in fact be made independent of it. It is also true that away from $\Gamma$ the spaces $J_{s} L^{2}(X)$ is the same as the space defined without including the lift of $H$ in the defining variety. That statement is non-trivial only near $D$.

The complications of the definitions are now compensated by the simplicity of the proof of the following

Theorem 3.8. - The spaces $J_{s} L_{\mathrm{loc}}^{2}(\tilde{X}, H) \cap L_{\mathrm{loc}}^{\infty}(\tilde{X})$ and $J_{s} L_{\mathrm{loc}}^{2}(X) \cap L_{\mathrm{loc}}^{\infty}(X)$ given by Definitions 3.5 and 3.6 respectively are $C^{\infty}$-modules and $C^{\infty}$-algebras.

In fact, we use the identification of the conormal spaces on the 'blown-up' side with $b$-Sobolev spaces (see Appendix B) and then apply the well known algebra properties of those spaces.

As in the earlier work on conormal regularity the difficult part is the propagation theorem. For the interior problem it is proved in Theorem 7, page 1026 of [34]:

Theorem 3.9. - If the variety $\mathfrak{S}_{5}$ in $\widetilde{X}_{5}$ is given by (3.9) with $\beta$ replaced by $\beta_{5}$, and

$$
P u=f \text { in } \tilde{X}, u \upharpoonright_{\tilde{X}_{-}}=0, f \in\left(\beta_{5}\right)_{*}\left(I_{k} L_{\nu_{5}}^{2}\left(\widetilde{X}_{5}, \varsigma_{5}\right)\right)
$$

then

$$
u, D u \in\left(\beta_{5}\right)_{*}\left(I_{k} L_{\nu_{5}, l o c}^{2}\left(\widetilde{X}_{5}, \varsigma_{5}\right)\right)
$$

The main result of this paper is the propagation theorem for the Dirichlet problem and the space $J_{k} L^{2}(X)$ - see chapter 8 . A simpler refinement of Theorem 3.9 to $J_{k} L^{2}(\widetilde{X}, H)$ will be given in chapter 5 .

## 4. MICROLOCALLY CHARACTERIZED SPACES OF DISTRIBUTIONS

4.1. The important notion of solutions to linear real principal type equations associated to Lagrangian submanifolds was introduced in [13]. That followed a rich tradition in geometric optics and semi-classical analysis already exploited in [17, 23] and generalized the notion of oscillatory solutions by recasting it in terms of propagation of singularities. The control of multiplicative properties required in the study of non-linear problems made it necessary to introduce a larger class of geometrically defined marked Lagrangian distributions [27, 30]. Additional motivation came also from the study of operators with double characteristics and of singular Radon transforms [12, 38]. For the basic material needed in this paper we refer the reader to the presentation in [34], chapter 9 , while a proper development of the theory will appear in [30].

The purpose of this section is to extend the notion of marked Lagrangian distributions in two directions: to sub- and super-marked Lagrangian spaces, the closely related spaces associated with Lagrangian manifolds with boundaries, and to marked Lagrangian spaces on a manifold with boundary, $M$. In the last case we allow, unlike in [34], certain Lagrangians which are not smooth in ${ }^{b} T^{*} M \backslash 0$. These two directions are rather independent at this point, with a rather special connection, however, which will be exploited in chapter 7 .

To make this section self-contained we start with the general discussion of marked spaces, conducted for simplicity in the case of infinite regularity. We shall then proceed with the more detailed theory of $\frac{1}{2}$ - and 2 -marked spaces. For our purposes it will be sufficient to consider only markings by a single submanifold.

Let $M$ be a $C^{\infty} n$-manifold without boundary and let $\Lambda_{0}, \Lambda_{1} \subset T^{*} M \backslash 0$ be conic Lagrangian $C^{\infty}$ submanifolds, intersecting cleanly. Then $K=\Lambda_{0} \cap \Lambda_{1}$ is a $C^{\infty}$ embedded conic submanifold of $\Lambda_{0}$, and any such $K$ can be obtained locally as a clean intersection with a Lagrangian submanifold $\Lambda_{1}$. We define the following marked Lagrangian varieties associated to this geometry:

$$
\begin{equation*}
\mathscr{L}_{0}=\Lambda_{0}, \quad \mathscr{L}_{1}=\left\{\Lambda_{0}, K\right\}, \quad \mathscr{L}_{\infty}=\Lambda_{0} \sqcup \Lambda_{1}=\left\{\Lambda_{0} \cup \Lambda_{1} \backslash \Lambda_{0} \cap \Lambda_{1}, \Lambda_{0} \cap \Lambda_{1}\right\} \tag{4.1}
\end{equation*}
$$

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and the corresponding $\Psi_{\mathrm{phg}}^{0}(M)$-modules of first order pseudodifferential operators:

$$
\begin{align*}
\Psi_{\mathrm{phg}}^{1}\left(M ; \mathfrak{L}_{t}\right)= & \left\{A \in \Psi_{\mathrm{phg}}^{\prime}(M):\left.\sigma_{1}(A)\right|_{L}=0, H_{\sigma_{1}}(A)\right.  \tag{4.2}\\
& \text { is tangent to } \left.L \text { for all } L \in \mathscr{L}_{t}\right\}, \quad t=0,1, \infty .
\end{align*}
$$

The $H_{(s)}(M)$-based spaces of (marked) Lagrangian distributions associated to the varieties (4.1) are defined as

$$
\begin{align*}
I H_{(s)}\left(M ; \mathscr{L}_{t}\right)= & \left\{u \in H_{(s)}(M): A_{1} \ldots A_{l} u \in H_{(s)}(M)\right. \\
& \text { for } \left.A_{i} \in \Psi_{\mathrm{phg}}^{1}\left(M ; \mathscr{L}_{t}\right) \text { and } l \in \mathbb{N}\right\}, \quad t=0,1, \infty . \tag{4.3}
\end{align*}
$$

The space on which the iterated regularity is based, $H_{(s)}$ can clearly be replaced by $H_{(s)}^{\text {comp }}$ or $H_{(s)}^{\text {loc }}$. Suppose now that $K=\Lambda_{0} \cap \Lambda_{1}$ is assumed to be an embedded hypersurface in $\Lambda_{0}$. Locally, we can use the following model for the geometry (which is a special case of a more general extension of Darboux's theorem [38]): $M \subset \mathbb{R}^{n}$ open, $0 \in M$ and

$$
\begin{equation*}
\Lambda_{0}=T_{0}^{*} \mathbb{R}^{n} \backslash 0, \quad \Lambda_{1}=N^{*}\left\{x: x_{1}=\cdots=x_{n-1}=0\right\} \tag{4.4}
\end{equation*}
$$

In this case one easily sees that $\Psi_{\mathrm{phg}}^{1}\left(M ; \mathscr{L}_{t}\right)$ is the $\Psi_{\mathrm{phg}}^{0}(M)$-span of

$$
\begin{gather*}
x_{i} D_{x_{j}}  \tag{4.5}\\
i, j=1, \ldots, n \quad t=0  \tag{4.6}\\
x_{i} D_{x_{j}}, x_{n} D_{x_{n}}, x_{n}^{2} D_{x_{j}}  \tag{4.7}\\
x_{i} D_{x_{j}}, x_{n} D_{x_{n}} \quad i=1, \ldots, n-1, j=1, \ldots, n \quad t=1 \\
\end{gather*}
$$

Thus

$$
\begin{equation*}
u \in I H_{(s)}\left(M ; \mathscr{L}_{t}\right) \Longleftrightarrow \widehat{u}^{\wedge} \in S L^{2}\left(m_{s}, g_{(1+t)^{-1}}\right) \tag{4.8}
\end{equation*}
$$

where $m_{s}=\left(1+|\xi|^{2}\right)^{\frac{s}{2}}$,

$$
g_{\alpha}=\frac{\left|d \xi^{\prime}\right|^{2}}{1+|\xi|^{2}}+\frac{d \xi_{n}^{2}}{\left(1+\left|\xi^{\prime}\right|^{2}\right)^{\alpha}+\xi_{n}^{2}+1}
$$

and $S L^{2}$ are the $L^{2}\left(\mathbb{R}^{n}\right)$ based symbols (replacing $L^{\infty}$ by $L^{2}$ in Sect. 18.4 of [14]). Since we have

$$
I H_{(s)}\left(M ; \mathfrak{L}_{0}\right) \subset I H_{(s)}\left(M, \mathfrak{\Re}_{1}\right) \subset I H_{(s)}\left(M, \mathscr{L}_{\infty}\right)
$$

one expects that changing of $\alpha=(1+t)^{-1}$ in (4.8) will lead to new classes of distributions which we will call sub-marked Lagrangian for $0<t<1$ and supermarked Lagrangian distributions for $1<t<\infty$. Since, unlike (4.2) and (4.3), the definitions using the model metrics $g_{\alpha}$ are not a priori invariant we might have to allow for more geometric information. This will indeed be the case for the super marked distributions.

To start we need to review the definition of the $k$-jet bundle of a manifold $Y$. Using the identification of $Y$ with the diagonal in $Y \times Y$,

$$
Y \simeq \Delta_{Y}=\{(y, y): y \in Y\} \subset Y \times Y
$$

we define for $k>0$ the following vector bundle over $Y$ :

$$
\begin{gather*}
J^{k}(Y)^{*}=\mathscr{J}(Y) / \mathscr{J}(Y)^{k+1}, \quad \mathscr{J}(Y) \subset C^{\infty}(Y \times Y)  \tag{4.9}\\
\mathscr{g}(Y)=\left\{f \in C^{\infty}(Y \times Y):\left.f\right|_{\Delta_{Y}}=0\right\} \tag{4.10}
\end{gather*}
$$

Then $J^{k}(Y)=\left(J^{k}(Y)^{*}\right)^{*} \longrightarrow Y$ is the $k$-jet bundle. If $Y \subset X$ is an embedded submanifold of a $C^{\infty}$ manifold $X$, then using the pull back of $\iota: Y \times Y \hookrightarrow X \times X$ we obtain a natural mapping $J^{k}(X)^{*} \longrightarrow J^{k}(Y)^{*}$ and consequently $J^{k}(Y) \longrightarrow J^{k}(X)$. Restricting $J^{k}(X)$ to $Y$ we obtain $J_{Y}^{k}(X)$ the jet of $X$ at $Y$, with the natural inclusion $J^{k}(Y) \subset J_{Y}^{k}(X)$. The inclusion map can be replaced by a more general $f: X_{1} \longrightarrow X_{2}$ which then induces

$$
J^{k} f: J^{k}\left(X_{1}\right) \longrightarrow J^{k}\left(X_{2}\right)
$$

If $X=X_{1}=X_{2}$ and $Y_{1} \subset Y \subset X$ are embedded submanifolds, we say that $f$ preserves the $k$-th jet of $Y$ at $Y_{1}$ if

$$
J^{k} f\left(J_{Y_{1}}^{k}(Y)\right) \subset J_{Y_{1}}^{k}(Y)
$$

where $J_{Y_{1}}^{k}(Y)$ is naturally included in $J_{Y_{1}}^{k}(X)$.
We now extend the notion of marked Lagrangian varieties of the type shown in (4.1) to the jet Lagrangian varieties:

$$
\begin{equation*}
\mathfrak{L}_{t}^{0}=\left\{\Lambda_{0}, J_{\Lambda_{0} \cap \Lambda_{1}}^{-([-t]+1)} \Lambda_{1}, t\right\} \tag{4.11}
\end{equation*}
$$

where we include the index $t, 0 \leq t<\infty$ in the variety and allow the convention $J^{0}(Y) \simeq Y$.

Definition 4.1. - For $0 \leq t<\infty$ and $\Lambda_{0}, \Lambda_{1}$ given by (4) we define the $t$-marked Lagrangian spaces associated to the jet Lagrangian varieties $\mathscr{L}_{t}^{0}$ given by (8) as

$$
I H_{(s)}\left(M, \mathfrak{L}_{t}^{0}\right)=\left\{u \in H_{(s)}(M):(\phi u)^{\wedge} \in S L^{2}\left(m_{s}, g_{(1+t)^{-1}}\right) \text { for } \phi \in C_{0}^{\infty}(M)\right\}
$$

For the harder finite regularity case and $t=2^{-1}, 1,2$, the invariance properties will be given in Proposition 4.1. In Proposition 4.3 we shall give, for the same $t$ 's and finite regularity, the proof of the following identity, quite easy in the case considered now

$$
\begin{equation*}
I H_{(s)}\left(M, \mathscr{L}_{\infty}\right)=I H_{(s)}\left(M, \mathscr{£}_{t}^{1}\right)+I H_{(s)}\left(M, \mathfrak{Q}_{\frac{1}{t}}^{0}\right), 0 \leq t \leq \infty \tag{4.12}
\end{equation*}
$$

where $\mathfrak{L}_{t}^{1}$ is the variety obtained from $\mathscr{L}_{t}$ in (8) by exchanging $\Lambda_{0}$ and $\Lambda_{1}$. This reflects the main rôle of the sub- and super-marked Lagrangian spaces, which is in providing decompositions of spaces associated to pairs of intersecting Lagrangians into terms which can be treated individually.

We conclude this introductory discussion by defining (marked) Lagrangian distributions associated to Lagrangians with boundaries, again only in the model case and with infinite order of regularity.

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Thus, in addition to the models given by (4.4) we consider

$$
\begin{equation*}
\Lambda_{0}^{ \pm}=\Lambda_{0} \cap\left\{ \pm \xi_{n} \geq 0\right\}, \quad \partial \Lambda_{0}^{ \pm}=\left\{\xi_{n}=0\right\} \cap \Lambda_{0}=\Lambda_{0} \cap \Lambda_{1} \tag{4.13}
\end{equation*}
$$

The marking will now indicate the microlocalization on one side of the boundary: the lower the marking the more localized the singularities are to $\Lambda_{0}^{ \pm}$. As we shall see in Proposition 4.2, for the case $t=2$, the invariance requires an additional marking in the transversal direction as can be seen from the uncertainty principle. As before, supermarking requires additional geometric information and we introduce the following jet Lagrangian varieties:

$$
\begin{equation*}
\varrho_{t}^{ \pm}=\left\{\Lambda_{0}^{ \pm}, t\right\}, t \geq 1 ; \varrho_{t}^{ \pm}=\left\{\Lambda_{0}^{ \pm}, J_{\partial \Lambda_{0}^{ \pm}}^{-\left(\left[-t^{-1}\right]+1\right)} \Lambda_{1}, t\right\}, \quad 0<t \leq 1 \tag{4.14}
\end{equation*}
$$

For $t=0$ we could consider $\Lambda_{0}^{ \pm} \sqcup \Lambda_{1}$ obtaining a slightly larger space than the one considered [38].

Definition 4.2. - For $\Lambda_{0}^{ \pm}$given by (4.13) and the jet varieties in (4.14) we define for $0<t<\infty$

$$
\begin{align*}
I H_{(s)}\left(M, \mathfrak{L}_{t}^{ \pm}\right)= & \left\{u \in I H_{(s)}\left(M, \mathfrak{L}_{t^{-1}}\right): \chi\left(\left|\xi_{1}\right|\right) \chi\left(-\left|\xi_{1}\right|^{-\frac{t}{1+t}} \xi_{n}\right)\left(\left|\xi_{1}\right|^{-\frac{t}{1+t}} \xi_{n}\right)^{N}(\phi u)^{\wedge}\right. \\
& \left.\in S L^{2}\left(m_{s}, g_{\left(\frac{t}{1+t}\right)}\right), \text { for } N \geq 0 \text { and } \phi \in C_{0}^{\infty}(M)\right\} \tag{4.15}
\end{align*}
$$

and where $\mathfrak{L}_{t}, m_{s}, g_{\alpha}$ are given by (4.14) and (4.8), $\chi \in C^{\infty}(\mathbb{R})$,

$$
\chi(t)= \begin{cases}0, & t<1  \tag{4.16}\\ 1, & t>2\end{cases}
$$

The decomposition (4.10) can now be refined to

$$
\begin{equation*}
I H_{(s)}\left(M, \mathscr{\Re}_{0}\right)=I H_{(s)}\left(M, \mathfrak{£}_{t}^{1}\right)+\sum_{ \pm} I H_{(s)}\left(M, \mathfrak{\varrho}_{t}^{ \pm}\right), 0<t<\infty . \tag{4.17}
\end{equation*}
$$

4.2. We shall now develop carefully the ideas presented in the first part of this section for the cases $t=\frac{1}{2}, 1,2$ and finite regularity. We recall that the marked Lagrangian distributions of [30] constitute the case $t=1$.

For the applications in chapter 7 it is convenient to state the definitions for a slightly different model:

$$
\begin{equation*}
\Lambda_{0}=N^{*}\left\{x_{1}=x_{2}=0\right\} \subset T^{*} \mathbb{R}^{n} \backslash 0, \quad \Lambda_{1}=N^{*}\left\{x_{1}=0\right\} \subset T^{*} \mathbb{R}^{n} \backslash 0 \tag{4.18}
\end{equation*}
$$

with the coordinates $x \in \mathbb{R}^{n}$ written as $\left(x_{1}, x_{2}, x^{\prime}\right), x^{\prime} \in \mathbb{R}^{n-2}$. For $M \subset \mathbb{R}^{n}$, open, $0 \in M$, we give

Definition 4.3. - Let $\Lambda_{0}$ and $\Lambda_{1}$ be given by (4.18). Then

$$
\begin{aligned}
& I_{2 k} H_{(s)}\left(M ; \Lambda_{0}, \Lambda_{0} \cap \Lambda_{1}, \frac{1}{2}\right)=\left\{u \in H_{(s)}(M) ; D_{x^{\prime}}^{\alpha^{\prime}}\left(x_{1} D_{x_{1}}\right)^{k_{1}}\left(x_{1} D_{x_{2}}\right)^{k_{2}}\left(x_{2} D_{x_{2}}\right)^{k_{3}} x_{2}^{k_{4}} u\right. \\
& \left.\in H_{\left(s+\frac{2}{3} k_{4}\right)}(M), \quad\left|\alpha^{\prime}\right|+k_{1}+k_{2}+k_{3}+\frac{2}{3} k_{4} \leq 2 k\right\} \\
& I_{2 k+1} H_{(s)}\left(M ; \Lambda_{0}, \Lambda_{0} \cap \Lambda_{1}, \frac{1}{2}\right)= \\
& \quad\left[I_{2(k+1)} H_{(s)}\left(M ; \Lambda_{0}, \Lambda_{0} \cap \Lambda_{1}, \frac{1}{2}\right), I_{2 k} H_{(s)}\left(M ; \Lambda_{0}, \Lambda_{0} \cap \Lambda_{1}, \frac{1}{2}\right)\right]_{\frac{1}{2}},
\end{aligned}
$$

and

$$
\begin{aligned}
I_{k} H_{(s)}\left(M ; \Lambda_{0}, J_{\Lambda_{0} \cap \Lambda_{1}}^{1}\right. & \left.\Lambda_{1}, 2\right)= \\
& \left\{u \in H_{(s)}(M): D_{x^{\prime}}^{\alpha^{\prime}}\left(x_{1} D_{x_{1}}\right)^{k_{1}}\left(x_{1} D_{x_{2}}\right)^{k_{2}}\left(x_{2} D_{x_{2}}\right)^{k_{3}} x_{2}^{k_{4}} u\right. \\
& \left.\in H_{\left(s+\frac{1}{3} k_{4}\right)}(M),\left|\alpha^{\prime}\right|+k_{1}+k_{2}+k_{3}+\frac{1}{3} k_{4} \leq k .\right\}
\end{aligned}
$$

One should observe that for $k=\infty$ these definitions coincide with Definition 4.1. In this case, as before, we could use the variety notation $\mathscr{L}_{t}, t=\frac{1}{2}, 2$. In both suband super-marked cases above, it is crucial that the operators appear with variable weights in the stability conditions. Thus we say that $\left\langle D_{x}\right\rangle^{\frac{2}{3}} x_{2}$ has weight $\frac{2}{3},\left\langle D_{x}\right\rangle^{\frac{1}{3}} x_{2}$, weight $\frac{1}{3}$. The correct filtration is based on operators of order 1 and that explains the need for fractional weights for operators of different orders.

We now want to verify the invariance of the $t$-marked spaces. Thus, let $\Gamma \subset T^{*} \mathbb{R}^{n} \backslash 0$ be a connected open conic neighbourhood of $m_{0}=(0 ;(1,0, \ldots, 0))$ and let us consider a canonical transformation

$$
\begin{equation*}
\chi: \Gamma \longrightarrow T^{*} \mathbb{R}^{n} \backslash 0, \chi\left(m_{0}\right)=m_{0} \tag{4.19}
\end{equation*}
$$

preserving the model geometry (4.18):

$$
\begin{equation*}
\chi\left(\Lambda_{0} \cap \Gamma\right) \subset \Lambda_{0}, J^{l} \chi\left(J_{\Lambda_{0} \cap \Lambda_{1}}^{l} \Lambda_{1} \cap J_{\Gamma}^{l}\left(T^{*} \mathbb{R}^{n} \backslash 0\right)\right) \subset J_{\Lambda_{0} \cap \Lambda_{1}}^{l} \Lambda_{1} \tag{4.20}
\end{equation*}
$$

For such $\chi$ we have
Proposition 4.4. - Let $F$ be a Fourier Integral Operator of order 0 associated to a canonical transformation $\chi$ satisfying (4.20) with $l=-([-t]+1)$. Then, if $t=\frac{1}{2}, 1,2$

$$
F: I_{k} H_{(s)}\left(M, \Lambda_{0}, J_{\Lambda_{0} \cap \Lambda_{1}}^{-([-t]+1)} \Lambda_{1}, t\right) \longrightarrow I_{k} H_{(s)}\left(M, \Lambda_{0}, J_{\Lambda_{0} \cap \Lambda_{1}}^{-[-t]+1} \Lambda_{1}, t\right)
$$

Proof. - We will give the proof in the case $t=2$. When $t=1$, the invariance is clear from the symplectically invariant definitions (4.2), (4.3) and the case $t=\frac{1}{2}$ is similar and simpler, as in particular it does not involve the jet bundle. It can also be derived by the methods of [30] or by using the calculus of Lagrangian distributions of class $I_{\frac{2}{3}}$ in [14], Sect. 25.1.

For $t=2$, the argument is easier in the model given by (4.4) and it is clear that by composition of Fourier Integral Operators we will obtain the invariance in the general case and in particular in the model given by (4.18). This means that, by Egorov's theorem, we change the operators in Definition 4.3 to

$$
\begin{align*}
& \left(x^{\prime}\right)^{\alpha^{\prime}} D_{x_{1}}^{\left|\alpha^{\prime}\right|}\left(x_{1} D_{x_{1}}\right)^{k_{1}}\left(x_{1} D_{x_{2}}\right)^{k_{2}}\left(x_{2} D_{x_{2}}\right)^{k_{3}} x_{2}^{k_{4}} u \in H_{\left(s+\frac{1}{3} k_{4}\right)}(M), \\
& \left|\alpha^{\prime}\right|+k_{1}+k_{2}+k_{3}+\frac{1}{3} k_{4} \leq k, \quad W F^{(k+s)}(u) \subset \Gamma . \tag{4.21}
\end{align*}
$$

Without the loss of generality, we can also assume that $s=0$. As in chapter 9 of [30] we first observe that if $\chi^{-1}(y, \eta)=(X(y, \eta), \Xi(y, \eta))$ then

$$
F u(y)=\frac{1}{(2 \pi)^{n}} \int e^{i\langle y, \eta\rangle} b(y, \eta) \hat{u}(\Xi(y, \eta)) d \eta, \quad b \in S^{0}\left(\mathbb{R}^{n} ; \mathbb{R}^{n}\right)
$$

and we start by specifying the properties of $\Xi=\left(\Xi_{1}, \Xi_{2}, \Xi^{\prime}\right), \Xi^{\prime} \in \mathbb{R}^{n-2}$, implied by the jet condition (4.20):

$$
\begin{align*}
& \Xi_{1}(y, \eta)=\eta_{1} h_{1}(y, \eta), h_{1}\left(m_{0}\right) \neq 0, h_{1} \text { homogeneous of degree } 0 \\
& \Xi_{2}(y, \eta)=\eta_{2} h_{2}(y, \eta)+y_{1} g_{1}(y, \eta)+y_{2}^{2} g_{2}(y, \eta)+\left\langle y^{\prime}, g(y, \eta)\right\rangle, \\
& h_{2}(y, \eta)=h_{20}(\eta)+y_{2} h_{22}(\eta, y), h_{2}\left(m_{0}\right) \neq 0, \tag{4.22}
\end{align*}
$$

$h_{i}$ homogeneous of degree $0, g_{i}, g$ homogeneous of degree 1 ,

$$
\left|\Xi^{\prime}(y, \eta)\right| \leq C\left|\eta_{1}\right| \text { for }(y, \eta) \in \widetilde{\Gamma}
$$

where $\widetilde{\Gamma} \subset \Gamma$ is another conic neighbourhood of $m_{0}, \operatorname{supp} b \subset \widetilde{\Gamma}$.
To illustrate the simple idea of the proof, let us consider the following integral:

$$
I\left(y_{1}, y_{2}\right)=\frac{1}{(2 \pi)^{2}} \int e^{i\langle y, \eta\rangle} a\left(\eta_{1}, \eta_{2}+y_{2}^{k} \eta_{1}\right) d \eta, \quad a \in S L^{2}\left(1, g_{\frac{1}{3}}\right)
$$

By taking the Fourier transform in the $y_{1}$ variable and making a substitution

$$
\lambda=\eta_{1}, y=\lambda^{\frac{1}{3}} y_{2}, \eta=\lambda^{\frac{1}{3}} \eta_{2}
$$

(that is, introducing a nonhomogeneous blow-up on the Fourier transform side, see chapter 7), we obtain

$$
I_{1}(\lambda, y)=\frac{1}{2 \pi} \int e^{i y \eta} a_{1}\left(\lambda, \eta+\lambda^{1-\frac{1}{3}(k+1)} y^{k}\right) \lambda^{\frac{1}{3}} d \eta
$$

with $a_{1}$ stable under $\lambda D_{\lambda}, \eta D_{\eta}, D_{\eta}$ with weights $1,1, \frac{1}{3}$ respectively. If $1-(k+1) / 3 \leq$ 0 , that is $k \geq 2$, one easily checks the stability of $I_{1}$ under the same operators and with the same weights.
Returning to the actual situation, we observe that in view of (4.21) it suffices to check that

$$
\begin{equation*}
\prod_{i \neq 2}\left(y_{i} D_{y_{1}}\right)^{k_{i}}\left(y_{2} D_{y_{2}}\right)^{k_{2}} y_{2}^{l} F u(y) \in H_{\left(\frac{1}{3} l\right)}(M) \tag{4.23}
\end{equation*}
$$

if

$$
\begin{equation*}
\prod_{i \neq 2}\left(\eta_{1} D_{\eta_{i}}\right)^{k_{i}}\left(\eta_{2} D_{\eta_{2}}\right)^{k_{2}}\left(\eta_{1}^{\frac{1}{3}} D_{\eta_{2}}\right)^{l} \chi\left(\eta_{1}\right) D_{\eta}^{\alpha} \hat{u}(\eta) \in L^{2}\left(\mathbb{R}^{n}\right), \quad u \in \mathfrak{G}^{\prime}(M) \tag{4.24}
\end{equation*}
$$

where in both cases $\Sigma k_{i}+\frac{1}{3} l \leq k$ and $\alpha$ is arbitrary. Using (4.22) we write

$$
\begin{align*}
F u(y)=\frac{1}{(2 \pi)^{n}} & \int e^{i\langle y, \eta\rangle} b(y, \eta) \times \\
& \hat{u}\left(\eta_{1} h_{1}, \eta_{2} h_{20}+\eta_{2} y_{2} h_{22}+y_{1} g_{1}+y_{2}^{2} g_{2}+\left\langle y^{\prime}, g\right\rangle, \Xi^{\prime}(y, \eta)\right) d \eta \tag{4.25}
\end{align*}
$$

where $h_{1} \neq 0$ and $h_{20} \neq 0$ in $\widetilde{\Gamma}$. Thus, we can replace $\hat{u}$ by $a\left(\eta_{1}, \eta_{2}+\eta_{2} y_{2} \bar{h}_{22}+\right.$ $y_{1} \bar{g}_{1}+y_{2}^{2} \bar{g}_{2}+\left\langle y^{\prime}, \bar{g}\right\rangle, \Xi^{\prime}(y, \eta)$ ), with $a(\eta)$ satisfying (4.24). We denote the integral obtained this way by $F(a, b)$ and we show that the application of the operators in (4.23) changes $F(a, b)$ to $F(\tilde{a}, \tilde{b})$ with $\tilde{b} \in S^{0}\left(\mathbb{R}^{n}, \mathbb{R}^{n}\right)$, or $S^{-\frac{1}{3}}\left(\mathbb{R}^{n}, \mathbb{R}^{n}\right) \operatorname{supp} \tilde{b} \subset \widetilde{\Gamma}$ and $\tilde{a}$ satisfying (4.24) with $k$ decreased depending on the weight of the operator applied (1 or $\frac{1}{3}$ ). We start with

$$
\begin{aligned}
& y_{2} F(a, b)(y)=\frac{1}{(2 \pi)^{n}} \int e^{i y \eta}\left(-D_{\eta_{2}}\right)(b a) d \eta \equiv \\
& \frac{-1}{(2 \pi)^{2}} \int e^{i\langle y, \eta\rangle}\left(1+y_{2} \bar{h}_{22}+\eta_{2} y_{2}\left(\bar{h}_{22}\right)_{\eta_{2}}^{\prime}+y_{1}\left(\bar{g}_{1}\right)_{\eta_{2}}^{\prime}+y_{2}^{2}\left(\bar{g}_{2}\right)_{\eta_{2}}^{\prime}+\left\langle y^{\prime},(\bar{g})_{\eta_{2}}^{\prime}\right\rangle\right)\left(D_{\eta_{2}} \tilde{a}\right) b d \eta
\end{aligned}
$$

where we omitted the terms $F\left(a, b_{1}\right), b_{1} \in S^{0}$. The term in brackets can be absorbed into $b$, while $a$ can be replaced by $\tilde{a}=\eta_{1}^{\frac{1}{3}} D_{\eta_{2}} a$. Thus we obtain $F(\tilde{a}, \tilde{b})$ with $\tilde{b} \in$ $S^{-\frac{1}{3}}\left(\mathbb{R}^{n}, \mathbb{R}^{n}\right)$ and $\tilde{a}$ satisfying (4.24) with $k$ replaced by $k-\frac{1}{3}$. Boundedness properties on Sobolev spaces for Fourier Integral Operators show that $F(\tilde{a}, \tilde{b}) \in H_{\left(s+\frac{1}{3}\right)}(M)$ if $F(a, b) \in H_{(s)}(M)$.
We now consider

$$
\begin{align*}
& \left(y_{2} D_{y_{2}}\right) F u=\frac{1}{(2 \pi)^{n}} \int e^{i y \eta} y_{2} \cdot \eta_{2}(b a) d \eta+\frac{1}{(2 \pi)^{n}} \int e^{i y \eta} y_{2} \times \\
& \quad\left[\left(\eta_{2} y_{2}\left(\bar{h}_{22}\right)_{y_{2}}^{\prime}+\eta_{2} \bar{h}_{22}+y_{1}\left(g_{1}\right)_{y_{2}}^{\prime}+\left\langle y^{\prime}, g_{y_{2}}\right\rangle\right)\left(D_{\eta_{2}} a\right)+\left\langle\Xi_{y_{2}}^{\prime}, D_{\eta^{\prime}}\right\rangle a\right] b d \eta \tag{4.26}
\end{align*}
$$

Since the needed estimate is local near $y_{1}=y_{2}=0$, we can, without the loss of generality, shrink the support of $b, \widetilde{\Gamma}$, so that $\left|y_{2} h_{22}(\eta)\right|<\frac{1}{2}$ in $\widetilde{\Gamma}$. Thus by changing $b$ to $\left(1+y_{2} h_{22}(\eta)\right)^{-1} b$, we can introduce the factor $\left(1+y_{2} h_{22}(\eta)\right)$ in the integrand on the right hand side of (4.26). Neglecting the terms of the form $F\left(a, b_{1}\right), b_{1} \in S^{0}$, the first term in the right hand side of (4.26) can then be written as

$$
\begin{align*}
& -\frac{1}{(2 \pi)^{n}} \int e^{i\langle y, \eta\rangle}\left(\eta_{2}+\eta_{2} y_{2} \bar{h}_{22}+y_{1} \bar{g}_{1}+y_{2}^{2} \bar{g}_{2}+\left\langle y^{\prime}, \bar{g}\right\rangle\right)\left(D_{\eta_{2}} a\right) \cdot b d \eta \\
& +\frac{1}{2 \pi)^{n}} \int e^{i\langle y, \eta\rangle}\left(y_{1} y_{2} \bar{g}_{1}+y_{2}^{3} \bar{g}_{2}+\left\langle y^{\prime}, \bar{g}\right\rangle\right) a b d \eta \tag{4.27}
\end{align*}
$$

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The first integral is of the form $F(\tilde{a}, b)$ with $\tilde{a}=\eta_{2} D_{\eta_{2}} a$ ( $\eta_{2}$ becomes the first term in brackets), which satisfies (4.24) with $k$ decreased by 1.
In the second term we integrate by parts using the stability under $\eta_{1} D_{\eta_{1}}, \eta_{1} D_{\eta_{2}}^{3}$ and $\eta_{1} D_{\eta^{\prime}}^{\alpha},|\alpha|=1$ in (4.24) for $y_{1}\left(y_{2} \bar{g}_{1}\right), y_{2}^{3}\left(\bar{g}_{2}\right)$ and $\left\langle y^{\prime}, \bar{g}\right\rangle$ respectively. In fact, let us verify this for $\eta_{1} D_{\eta_{1}}$ :

$$
\begin{align*}
& \left(\eta_{1} D_{\eta_{1}}\right)\left[a\left(\eta_{1}, \eta_{2}+\eta_{2} y_{2} h_{22}+y_{1} g_{1}+y_{2}^{2} g_{2}+\left\langle y^{\prime}, g\right\rangle, \Xi^{\prime}\right)\right] \equiv  \tag{4.28}\\
& {\left[y_{2} \eta_{2}\left(\eta_{1}\left(\bar{h}_{22}\right)_{\eta_{1}}^{\prime}\right)+\eta_{1} y_{1}\left(\bar{g}_{1}\right)_{\eta_{1}}^{\prime}+\eta_{1} y_{2}^{2}\left(\bar{g}_{2}\right)_{\eta_{1}}^{\prime}+\left\langle y^{\prime}, \bar{g}_{\eta_{1}}^{\prime}\right\rangle\right] D_{\eta_{2}} a}
\end{align*}
$$

where we omitted the terms $\eta_{1} D_{\eta_{1}} a$ and $\eta_{1} D_{\eta^{\prime}} a$ as then we can use the operators in (4.24). Since we have the stability under $\eta_{1} D_{\eta_{2}}^{3}$ with weight 1 in (4.24), the goal is to bring the number of $\eta_{2}$ derivatives falling on $a$ to three. For the first term in (4.28) we use (4.27) again, with $b$ in the first term replaced by $y_{2}\left(\eta_{1}\left(h_{22}\right)_{\eta_{1}}^{\prime}\right) b \in S^{0}(\Gamma)$ and $a$ in the second term by $y_{2} D_{\eta_{2}} a$. For the remaining terms in (4.28), we use $y_{i} \exp (i\langle y, \eta\rangle)=D_{\eta_{i}} \exp (i\langle y, \eta\rangle)$ and integration by parts. In each case there is a gain in the number of $\eta_{2}$ derivatives, so after at most three applications we obtain $\eta_{1} D_{\eta_{2}}^{3} a$.
Since $\eta_{1}^{-1} g_{i}, \eta_{1}^{-1} \bar{g} \in S^{0}(\Gamma)$, the second integral is again of the form $F(\tilde{a}, b)$, with $b \in S^{0}$ and $\tilde{a}$ satisfying (4.24) with $k$ decreased by 1 . The analysis of the second term on the right hand side of (4.26) and the verification of the stability under the remaining operators in (4.23) are similar and are left to the reader.

We shall now present the finite regularity analogue of Definition 4.2. In that we restrict ourselves to the case relevant in our applications, $t=2$.

Definition 4.5. - Let $\Lambda_{0}$ and $\Lambda_{1}$ be given by (4.18) and let $M \subset \mathbb{R}^{n}$ be a bounded open set, $0 \in M$. If $\Lambda_{0}^{ \pm}=\Lambda_{0} \cap\left\{ \pm \xi_{n} \geq 0\right\}$, then

$$
\begin{aligned}
I_{2 k} H_{(s)}\left(M, \Lambda_{0}^{ \pm}, 2\right)= & \left\{u \in H_{(s)}(M): D_{x^{\prime}}^{\alpha^{\prime}}\left(x_{1} D_{x_{1}}\right)^{k_{1}}\left(x_{1} D_{x_{2}}\right)^{k_{2}}\left(x_{2} D_{x_{2}}\right)^{k_{3}} x_{2}^{k_{4}}\right. \\
& \left(\chi\left(\left|\xi_{1}\right|\right) \chi\left(\mp\left|\xi_{1}\right|^{-\frac{2}{3}} \xi_{2}\right)\left(\left|\xi_{1}\right|^{-\frac{2}{3}} \xi_{2}\right)^{k_{0}}(\phi u)^{\wedge}\right)^{\vee} \in H_{\left(s+\frac{2}{3} k_{4}\right)}(M) \\
& \text { for } \left.\phi \in C_{0}^{\infty}(M),\left|\alpha^{\prime}\right|+k_{1}+k_{2}+k_{3}+\frac{2}{3} k_{4}+\frac{1}{3} k_{0} \leq 2 k\right\}
\end{aligned}
$$

with $\chi \in C^{\infty}(\mathbb{R})$ satisfying (4.16) and $(\bullet)^{\vee}$ denoting the inverse Fourier transform. For odd orders of regularity $2 k+1, I_{2 k+1} H_{(s)}\left(M, \Lambda_{0}^{ \pm}, 2\right)$ is defined by complex interpolation as in Definition 4.3.

To study the invariance properties, let $\Gamma$ and $\chi$ be as in (4.19) with (4.20) replaced by

$$
\begin{equation*}
\chi\left(\Lambda_{0}^{ \pm} \cap \Gamma\right) \subset \Lambda_{0}^{ \pm} \tag{4.29}
\end{equation*}
$$

The analogue of Proposition 4.1 is:
Proposition 4.6. - Let $F$ be a Fourier Integral Operator of order 0 associated to the canonical transformation $\chi$ which satisfies (4.29). Then

$$
F: I_{k} H_{(s)}\left(M, \Lambda_{0}^{ \pm}, 2\right) \longrightarrow I_{k} H_{(s)}\left(M, \Lambda_{0}^{ \pm}, 2\right)
$$

Proof. - We can consider only $\Lambda_{0}^{+}$and assume that $k$ is even as the odd order case will follow by complex interpolation. We shall use the same model (4.4) as in the proof of Proposition 4.1 with $\Lambda_{0}^{ \pm}$given by (4.13) As in (4.21), Egorov's theorem provides the defining operators (see (4.34) below) and we can assume that

$$
W F^{(k)}(u) \subset \widetilde{\Gamma}, W F^{(k)}(F u) \subset \tilde{\Gamma}
$$

where $\widetilde{\Gamma} \subset \Gamma$ is a small conic neighbourhood of $m_{0}$. It is convenient however to use the representation of $F$ involving the generating function of the canonical transformation:

$$
\begin{align*}
& (y, \eta)=\chi(x, \xi) \Longleftrightarrow(y, \eta ; x, \xi)=\left(y, \phi_{y}^{\prime} ; \phi_{\xi}^{\prime}, \xi\right) \\
& F u(y)=\frac{1}{(2 \pi)^{n}} \int e^{i \phi(y, \xi)} b(y, \xi) \hat{u}(\xi) d \xi \tag{4.30}
\end{align*}
$$

The assumption (4.29) implies that

$$
\begin{equation*}
\phi(y, \xi)=y_{1} \xi_{1} h_{1}(y, \xi)+y_{2} \xi_{2} h_{2}(y, \xi)+\left\langle y^{\prime}, g(y, \xi)\right\rangle \tag{4.31}
\end{equation*}
$$

where $h_{1}, h_{2}$ are homogeneous of degree 0 and positive in $\Gamma, g$ is homogeneous of degree 1 with $g_{\xi^{\prime}}^{\prime}$ of rank $n-2$ in $\Gamma$.
We shall, as in the proof of Proposition 4.1, consider more general integrals, starting with the Fourier transform of (4.30):

$$
\begin{equation*}
G(a, b)(\eta)=\frac{1}{(2 \pi)^{n}} \int e^{i \phi(y, \xi)-i\langle y, \xi\rangle} \psi\left(\xi_{1}^{-1} \eta_{1}\right) b(y ; \xi, \eta) a(\xi) d \xi d y \tag{4.32}
\end{equation*}
$$

where $b \in S_{\frac{2}{3}, 0}^{-l}\left(\mathbb{R}^{n}, \mathbb{R}^{2 n}\right)$, supp $b \subset \Gamma_{1}$ with $\Gamma_{1}$ a small conic neighbourhood of $(0 ;(1,0, \ldots, 0),(1,0, \ldots 0))$ (with respect to the $\mathbb{R}_{+}$-action in the last $2 n$ coordinates). In fact $b$ is assumed to satisfy a stronger estimate:

$$
\begin{align*}
& \left|D_{\xi}^{\alpha} D_{\eta}^{\beta} b\right| \leq C_{\alpha \beta}\left(1+\xi_{1}+\eta_{1}\right)^{-\left(\left|\alpha^{\prime \prime}\right|+\left|\beta^{\prime \prime}\right|+\frac{2}{3}\left(\alpha_{2}+\beta_{2}\right)\right)-l}  \tag{4.33}\\
& l=l(b), \alpha^{\prime \prime}=\left(\alpha_{1}, \alpha_{3}, \ldots, \alpha_{n}\right), \beta^{\prime \prime}=\left(\beta_{1}, \beta_{3}, \ldots, \beta_{n}\right)
\end{align*}
$$

The cut-off $\psi \in C_{0}^{\infty}(\mathbb{R} \backslash 0)$, is chosen so that $\psi(t)=1$ for $\frac{1}{C}<t<C$ for some $C$. We assume also that $a$ satisfies

$$
\begin{align*}
& \chi\left(-\xi_{1}^{-\frac{2}{3}} \xi_{2}\right)\left(\xi_{1}^{-\frac{2}{3}} \xi_{2}\right)^{k_{0}}\left(\xi_{1} D_{\xi_{1}}\right)^{k_{1}} \xi_{1}^{\left|\alpha^{\prime}\right|} D_{\xi^{\prime}}^{\alpha^{\prime}}\left(\xi_{2} D_{\xi_{2}}\right)^{k_{2}}\left(\xi_{1}^{\frac{2}{3}} D_{\xi_{2}}\right)^{k_{3}} a \\
& \quad \in\left(1+\left|\xi_{1}\right|\right)^{\max (0,-p)} L^{2}\left(\mathbb{R}^{n}\right)  \tag{4.34}\\
& \frac{1}{3} k_{0}+k_{1}+k_{2}+\frac{2}{3} k_{3}+\left|\alpha^{\prime}\right| \leq \max (0, p), \quad p=p(a) \in \frac{1}{3} \mathbb{Z}
\end{align*}
$$

We observe that the boundedness of Fourier Integral Operators in the class $I_{\frac{2}{3}}^{-l}$, implies that

$$
\begin{equation*}
G(a, b) \in H_{(l(b)+\max (0,-p(a)))}(M) \tag{4.35}
\end{equation*}
$$

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The cut-off $\psi\left(\xi_{1}^{-1} \eta_{1}\right)$ in (4.32) can be inserted, at the expense of smooth error terms, in the integrand on the right hand side of

$$
\begin{equation*}
\widehat{F u}(\eta)=\frac{1}{(2 \pi)^{n}} \int e^{i \phi(y, \xi)-i\langle y, \eta\rangle} b(y, \xi) \hat{u}(\xi) d \xi d y \tag{4.36}
\end{equation*}
$$

since

$$
\begin{equation*}
\int e^{i \phi(y, \xi)-i\langle y, \eta\rangle} b(y, \xi)\left(1-\psi\left(\xi_{1}^{-1} \eta_{1}\right)\right) d y=\mathcal{O}\left(\left(1+\left|\xi_{1}\right|+\left|\eta_{1}\right|\right)^{-N}\right) \tag{4.37}
\end{equation*}
$$

for any $N>0$. In fact since $h_{1}(y, \xi)>0$ in $\widetilde{\Gamma}, h_{1}(0,(1,0, \ldots, 0))=1$, we have $(\phi(y, \xi)-\langle y, \eta\rangle)_{y_{1}}^{\prime}>c\left(1+\left|\xi_{1}\right|+\left|\eta_{1}\right|\right)$ in the support of $(1-\psi)\left(\xi_{1}^{-1} \eta_{1}\right)$ if $C$ and $c$ are appropriately chosen. Thus, standard oscillatory integral estimates (see Theorem 7.7 .1 of [14], as applied in Sect. 25.1 there) give (4.22). We conclude that $\widehat{F u}$ is essentially of the form (4.32) with $p(a)=k$, and we want to show that it also satisfies the estimate (4.34) with $p=k$.
If we apply the first operator in (4.34) to (4.32) we obtain

$$
\begin{equation*}
\chi\left(-\eta_{1}^{-\frac{2}{3}} \eta_{2}\right)\left(\eta_{1}^{-\frac{2}{3}} \eta_{2}\right) G(a, b)=I_{1}+I_{2} \tag{4.38}
\end{equation*}
$$

where the decomposition was obtained by inserting $1=\chi\left(\xi_{1}^{-\frac{2}{3}} \xi_{2}\right)+(1-\chi)\left(\xi_{1}^{-\frac{2}{3}} \xi_{2}\right)$ in the integrand:

$$
\begin{equation*}
I_{1}=\frac{1}{(2 \pi)^{n}} \int e^{i \phi(y, \xi)-i\langle y, \eta\rangle} \eta_{1}^{-\frac{2}{3}} \eta_{2} b(y ; \xi, \eta) \chi\left(-\eta_{1}^{-\frac{1}{3}} \eta_{2}\right) \chi\left(\xi_{1}^{-\frac{1}{3}} \xi_{2}\right) a(\xi) d \xi d y \tag{4.39}
\end{equation*}
$$

We shall now integrate by parts to put $I_{1}$ in the form (4.32) with new $a$ and $b$. Strictly speaking, we should also introduce a cutoff in $\xi^{\prime} / \xi_{1}$ similar to $\psi$, reducing the integration in $\xi^{\prime} / \xi_{1}$ to a compact set. For this we observe that since $g_{\xi^{\prime}}^{\prime}$ in (4.31) has rank $n-2, y^{\prime}$ can be expressed in terms of $\left\langle y^{\prime}, g_{\xi^{\prime}}^{\prime}\right\rangle \in \mathbb{R}^{n-2}$ and there exist differential operators $Q_{1}\left(y, \xi, D_{\xi^{\prime}}\right)$ and $Q_{2}\left(y, \xi, D_{\xi^{\prime}}\right)$ of order 1 with coefficients homogeneous of degree 1 in $\xi$ such that

$$
\begin{align*}
& e^{i \phi(y, \xi)-i\langle y, \eta\rangle}=\left(\eta_{2}-\xi_{2}\left(h_{2}+y_{2}\left(h_{2}\right)_{y_{2}}^{\prime}\right)\right)^{-1} \\
& \quad\left[\left(-D_{y_{2}}+Q_{2}\left(y, \xi, D_{\xi^{\prime}}\right)\right)-\left(h_{1}\right)_{y_{2}}^{\prime}\left(1-\left(\xi_{2} / \xi_{1}\right) \xi_{1}^{2}\left(h_{2}\right)_{\xi_{1}}^{\prime}\left(h_{1}\right)_{\xi_{2}}^{\prime}\right)^{-1} \times\right. \\
& \left.\quad\left(\xi_{1} D_{\xi_{1}}-\xi_{1}\left(h_{2}\right)_{\xi_{1}}^{\prime} \xi_{2} D_{\xi_{2}}+Q_{1}\left(y, \xi, D_{\xi^{\prime}}\right)\right)\right] e^{i \phi(y, \xi)-i\langle y, \eta\rangle} \tag{4.40}
\end{align*}
$$

The last term is obtained by writing

$$
\begin{align*}
& y_{1} e^{i \phi(y, \xi)-i\langle y, \eta\rangle}=\left(D_{\xi_{1}}-y_{2} \xi_{2}\left(h_{2}\right)_{\xi_{1}}^{\prime}-\left\langle y^{\prime}, \bar{g}_{\xi_{1}}^{\prime}\right\rangle\right) e^{i \phi(y, \xi)-i\langle y, \eta\rangle}=  \tag{4.41}\\
& {\left[D_{\xi_{1}}-\xi_{2}\left(h_{2}\right)_{\xi_{2}}^{\prime}\left(D_{\xi_{2}}-\xi_{1} y_{1}\left(h_{1}\right)_{\xi_{2}}-\left\langle y^{\prime}, \bar{g}_{\xi_{2}}^{\prime}\right\rangle\right)-\left\langle y^{\prime}, \bar{g}_{\xi_{1}}^{\prime}\right\rangle\right] e^{i \phi(y, \xi)-i\langle y, \eta\rangle}}
\end{align*}
$$

and observing that if $\widetilde{\Gamma}$ is small enough then $\left|\left(\xi_{2} / \xi_{1}\right) \xi_{1}^{2}\left(h_{2}\right)_{\xi_{1}}^{\prime}\left(h_{1}\right)_{\xi_{2}}^{\prime}\right|<\frac{1}{2}$.
Since $\left|\eta_{2}-\xi_{2}\left(h_{2}+y_{2}\left(h_{2}\right)_{y_{2}}^{\prime}\right)\right|>c \max \left\{\xi_{1}^{\frac{2}{3}},\left|\xi_{2}\right|+\left|\eta_{2}\right|\right\}$ in the support of the integrand, we obtain $I_{1}=G\left(a_{1}, b_{1}\right)$ where $b_{1} \in S_{\frac{2}{3}, 0}^{-l_{1}}\left(\mathbb{R}^{n} ; \mathbb{R}^{2 n}\right), l_{1}=l(b)+\frac{2}{3}$ (with the stronger
estimate (4.33) valid) and $a_{1}$ satisfies (4.34) with $p\left(a_{1}\right)=p(a)-1$. Note that the gain of regularity in $b_{1}$, compensates the decrease of $p$, so that the weight of the operator on $a$ is $\frac{1}{3}$, same as that of $\chi\left(-\eta_{1}^{-\frac{2}{3}} \eta_{2}\right) \eta_{1}^{-\frac{2}{3}} \eta_{2}$ - see (4.34) and (4.35).
The analysis of $I_{2}$ is similar but we now invoke the assumption that $a$ is stable under

$$
\begin{equation*}
(1-\chi)\left(\xi_{1}^{-\frac{2}{3}} \xi_{2}\right) \xi_{1}^{-\frac{2}{3}} \xi_{2} \tag{4.42}
\end{equation*}
$$

with weight $\frac{1}{3}$. To do that we again use (4.40) to write

$$
\begin{aligned}
I_{2}= & \frac{1}{(2 \pi)^{n}} \int\left(\xi_{1} / \eta_{1}\right)^{\frac{2}{3}} \xi_{1}^{-\frac{2}{3}}\left[\xi_{2}\left(h_{2}+y_{2}\left(h_{2}\right)_{y_{2}}^{\prime}\right)-D_{y_{2}}+Q_{2}\left(y, \xi, D_{\xi^{\prime}}\right)-\right. \\
& \left(h_{1}\right)_{y_{2}}^{\prime}\left(1-\left(\xi_{2} / \xi_{1}\right) \xi_{1}^{2}\left(h_{2}\right)_{\xi_{1}}^{\prime}\left(h_{1}\right)_{\xi_{2}}^{\prime}\right)^{-1}\left(\xi_{1} D_{\xi_{1}}-\xi_{1}\left(h_{2}\right)_{\xi_{1}}^{\prime} \xi_{2} D_{\xi_{2}}+\right. \\
& \left.\left.Q_{1}\left(y, \xi, D_{\xi^{\prime}}\right)\right)\right] e^{i \phi(y, \xi)-i\langle y, \eta\rangle} b(y ; \xi, \eta)(1-\chi)\left(\xi_{1}^{-\frac{2}{3}} \xi_{2}\right) \chi\left(-\eta_{1}^{\frac{2}{3}} \eta_{2}\right) a(\xi) d \xi d y
\end{aligned}
$$

and then integrate by parts. Thus $I_{2}=G\left(a_{2}, b_{2}\right)+G\left(a_{3}, b_{3}\right)$, where $a_{2}$ satisfies (4.34) with $p\left(a_{2}\right)=p(a)-\frac{1}{3}, b_{2}$ satisfies (4.33) with $l\left(b_{2}\right)=l(b)$ and $a_{3}$ satisfies (4.34) with $p\left(a_{3}\right)=p(a)-1, b_{3}$ satisfies (4.33) with $l\left(b_{3}\right)=l(b)+\frac{2}{3}$. In both cases the effective weights of operators are preserved.

The action of the remaining operators is similar and simpler. Thus we obtain

$$
\begin{align*}
& \chi\left(-\eta_{1}^{-\frac{2}{3}} \eta_{2}\right)\left(-\eta_{1}^{-\frac{2}{3}} \eta_{2}\right)^{k_{0}}\left(\eta_{1} D_{\xi_{1}}\right)^{k_{1}} \eta_{1}^{\left|\alpha^{\prime}\right|} D_{\eta^{\prime}}^{\alpha^{\prime}}\left(\eta_{2} D_{\eta_{2}}\right)^{k_{2}}\left(\eta_{1}^{\frac{2}{3}} D_{\eta_{2}}\right)^{k_{3}} G(a, b) \\
& =\sum_{i=1}^{N} G\left(a_{i}, b_{i}\right), \quad \frac{1}{3} k_{0}+k_{1}+k_{2}+\frac{2}{3} k_{3}+\left|\alpha^{\prime}\right| \leq p(a) \tag{4.43}
\end{align*}
$$

with $a_{i}$ satisfying (4.34) with $p\left(a_{i}\right) \leq p(a)$ and $b_{i} \in S_{\frac{2}{3}, 0}^{-l}\left(\mathbb{R}^{n} ; \mathbb{R}^{2 n}\right), l=\max (0,-p(a))$. Using (4.35) we conclude that $\widehat{F u}$ satisfies (4.34) with $p=k$.

Having established the invariance we can now give the general definition. Let $M$ be a $C^{\infty}$ manifold of dimension $n$ and let $\Lambda, \widetilde{\Lambda} \subset T^{*} M \backslash 0$ be conic Lagrangian submanifolds intersecting cleanly at $\Lambda \cap \widetilde{\Lambda}$, a hypersurface in $\Lambda$. For $t=\frac{1}{2}, 1,2$ we define the jet Lagrangian variety

$$
\begin{equation*}
\mathscr{L}_{t}=\left\{\Lambda, J_{\Lambda \cap \widetilde{\Lambda}}^{-([-t]+1)} \widetilde{\Lambda}, t\right\} . \tag{4.44}
\end{equation*}
$$

We also consider $\Lambda^{+} \subset \Lambda \subset T^{*} M \backslash 0$, a Lagrangian submanifold with boundary, $\partial \Lambda^{+}=\Lambda \cap \widetilde{\Lambda}$ and the variety

$$
\begin{equation*}
\mathfrak{L}^{+}=\left\{\Lambda^{+}, 2\right\} . \tag{4.45}
\end{equation*}
$$

We should remark that subsequently we may use either the notation $\mathfrak{L}_{t}$ or the explicit description of the variety as given by the right hand side of (4.44).

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Definition 4.7. - For $\Lambda, \widetilde{\Lambda}, \mathscr{L}_{t}$ above and $t=\frac{1}{2}, 1,2$ we define the space

$$
I_{k} H_{(s)}\left(M ; \mathscr{Q}_{t}\right)
$$

as consisting of $u \in \mathscr{D}^{\prime}(M)$ such that there exists $u_{0} \in I_{k} H_{(s)}(M ; \Lambda \backslash \Lambda \cap \widetilde{\Lambda})$, a countable covering of $\Lambda \cap \widetilde{\Lambda}$ by parametrizations

$$
\begin{gather*}
\chi_{j}: \Gamma_{j} \longrightarrow T^{*} \mathbb{R}^{n} \backslash 0, \quad \Lambda \cap \tilde{\Lambda} \subset \bigcup_{j} \Gamma_{j},  \tag{4.46}\\
\chi_{j}\left(\Gamma_{j} \cap \Lambda\right) \subset \Lambda_{0}, J^{-([-t]+1)} \chi_{j}: J_{\Lambda \cap \widetilde{\Lambda}}^{-([-t]+1)} \widetilde{\Lambda} \longrightarrow J_{\Lambda_{0} \cap \Lambda_{1}}^{-([-t]+1)} \Lambda_{1},
\end{gather*}
$$

where $\Lambda_{0}$ and $\Lambda_{1}$ are given by (4.18),

$$
F_{j} \in I^{0}\left(\mathbb{R}^{n}, M ;\left(\chi_{j}^{-1}\right)^{\prime}\right) \quad \text { and } \quad v_{j} \in I_{k} H_{(s)}\left(\mathbb{R}^{n} ; \Lambda_{0}, J_{\Lambda \cap \Lambda_{1}}^{-([-t]+1)} \Lambda_{1}, t\right)
$$

where the last space was given in Definition 4.3. The distributions $F_{j} v_{j}$ are assumed to have locally finite supports such that

$$
u-u_{0}-\sum_{j} F_{j} v_{j} \in H_{(k)}(M)
$$

If $\mathfrak{L}^{+}$is given by (4.45), the space

$$
I_{k} H_{(s)}\left(M ; \mathfrak{Q}^{+}\right)
$$

is defined in a similar way using $I_{k} H_{(s)}\left(\mathbb{R}^{n} ; \Lambda_{0}^{+}, 2\right)$ given in Definition 4.5.
We now proceed with the finite regularity case of the decompositions (4.12) and (4.17). For the marked Lagrangian distributions $(t=1)$ the proposition below was already established in [30]. Its proof here illustrates, in a computationally simple case, the general philosophy of relating microlocal and conormal spaces (see [25, 34, 42]), on which we shall rely heavily in chapter 6 and 7 . For the variety $\Lambda \sqcup \widetilde{\Lambda}=\{\Lambda \cup \widetilde{\Lambda} \backslash \Lambda \cap \widetilde{\Lambda}, \Lambda \cap \widetilde{\Lambda}\}$ we define $I_{k} H_{(s)}(M, \Lambda \sqcup \widetilde{\Lambda})$ by (4.3) with the obvious finite order modification.

Proposition 4.8. - Let $\Lambda$ and $\widetilde{\Lambda}$ be as in Definition 4.7. Then for $t=\frac{1}{2}, 1,2$ we have

$$
\begin{equation*}
I_{k} L_{\mathrm{loc}}^{2}(M ; \Lambda \sqcup \widetilde{\Lambda})=I_{k} L_{\mathrm{loc}}^{2}\left(M ; \mathfrak{L}_{t}\right)+I_{k} L_{\mathrm{loc}}^{2}\left(M ; \widetilde{\mathfrak{L}}_{\frac{1}{t}}\right) \tag{4.47}
\end{equation*}
$$

where $\mathscr{L}_{t}$ is given by (4.24) and $\widetilde{\mathscr{L}}_{t}$ is obtained by exchanging $\Lambda$ and $\widetilde{\Lambda}$ there. In addition, if $\Lambda=\cup \Lambda_{ \pm}, \Lambda_{+} \cap \Lambda_{-}=\partial \Lambda^{ \pm}=\Lambda \cap \tilde{\Lambda}$, then

$$
\begin{equation*}
I_{k} L_{\mathrm{loc}}^{2}(M, \Lambda \sqcup \tilde{\Lambda})=I_{k} L_{\mathrm{loc}}^{2}\left(M ; \tilde{\mathfrak{~}}_{2}\right)+\sum_{ \pm} I_{k} L_{\mathrm{loc}}^{2}\left(M ; \mathfrak{\varrho}^{ \pm}\right) \tag{4.48}
\end{equation*}
$$

with $\mathfrak{£}^{ \pm}$given by (4.45).

In the proof we shall only consider the case $t=\frac{1}{2}$ as $t=2$ is symmetric to it, while $t=1$ is discussed in Remark 4.12. The proof is based on a lifting of the right hand side of (4.47) to a conormal space. As the result is local and in view of the invariance, we take $M$ to be a neighbourhood of 0 in $\mathbb{R}^{n}$ and $\Lambda=\Lambda_{1}, \widetilde{\Lambda}=\Lambda_{0}$ given by (4.18).

With $M=\mathbb{R}^{n}$ we consider the following successive blow-ups of $\left\{x_{1}=x_{2}=0\right\}$

$$
\begin{equation*}
M_{2} \xrightarrow{\beta_{1-2}} M_{1} \xrightarrow{\beta_{3-1}} M, \quad \beta=\beta_{1-2} \circ \beta_{3-1}, \tag{4.49}
\end{equation*}
$$

where $M_{1}$ is a manifold with boundary and $M_{2}$ a manifold with corners:
$M_{1}=M \backslash\{0\} \sqcup\left(\mathbb{S}_{1-3}^{1} \times \mathbb{R}^{n-2}\right) \simeq \mathbb{R}_{+} \times \mathbb{S}_{1-3}^{1} \times \mathbb{R}^{n-2}, \mathbb{S}_{1-3}^{1}=\left\{\omega \in \mathbb{R}^{2}: \omega_{1}^{4}+\omega_{2}^{12}=1\right\}$, with the $C^{\infty}$ structure given by the second identification. The blow-down map $\beta_{3-1}: M_{1} \longrightarrow M$ is given by

$$
\begin{equation*}
\beta_{3-1}:\left(r, \omega, x^{\prime}\right) \longmapsto\left(r^{3} \omega_{1}, r \omega_{2}, x^{\prime}\right) \tag{4.50}
\end{equation*}
$$

The second resolution $M_{2}$ is obtained by the blow-up of $\beta_{3-1}^{*}\left\{x_{2}=0\right\} \cap \partial M_{1}$ :

$$
M_{2}=M_{1} \backslash\left(\beta_{3-1}^{*}\left\{x_{2}=0\right\} \cap \partial M_{1}\right) \sqcup\left(\mathbb{S}_{1_{-}-2}^{1} \times \mathbb{R}^{n-2}\right) \sqcup\left(\mathbb{S}_{1_{+}-2}^{1} \times \mathbb{R}^{n-2}\right)
$$

where now $\mathbb{S}_{1_{ \pm}-2}^{1}=\left\{\left(\theta_{1}, \theta_{2}\right): \pm \theta_{1} \geq 0, \theta_{1}^{4}+\theta_{2}^{2}=1\right\}$ (see Fig. 4.1) with the usual $C^{\infty}$ structure (see [28]). The intermediate blow-down map $\beta_{1-2}: M_{2} \longrightarrow M_{1}$ is given by $\beta_{1-2}\left(\rho, \theta, x^{\prime}\right)=\left( \pm \rho \theta_{1}, \rho^{2} \theta_{2}, x^{\prime}\right)$, where $\left(\rho, \theta, x^{\prime}\right) \in \mathbb{R}_{+} \times \mathbb{S}_{1_{+}-2}^{1} \times \mathbb{R}^{n-2}$, with coordinates near $\beta_{3-1}^{*}\left\{x_{2}=0, \pm x_{1}>0\right\} \cap \partial M_{1},\left(r, X_{2}, x^{\prime}\right)$ chosen so that $\beta_{3-1}\left(r, X_{2}, x^{\prime}\right)=\left( \pm r^{3}, r X_{2}, x^{\prime}\right)$.

Using the definition (4.3) we easily see that $u \in I_{k} L_{\text {loc }}^{2}\left(M ; \Lambda_{0} \sqcup \Lambda_{1}\right)$ if and only if

$$
\left(x_{1} D_{x_{1}}\right)^{k_{1}}\left(x_{2} D_{x_{2}}\right)^{k_{2}}\left(x_{1} D_{x_{2}}\right)^{k_{3}} D_{x^{\prime}}^{\alpha^{\prime}} u \in L_{\mathrm{loc}}^{2}(M), \quad \Sigma k_{i}+\left|\alpha^{\prime}\right| \leq k .
$$

This condition, Definition 4.3 of $I_{k} L_{\mathrm{loc}}^{2}\left(M ; \Lambda_{0}, J_{\Lambda_{0} \cap \Lambda_{1}}^{1}, 2\right)$ and the lifting of the vector fields in projective coordinates, give

Lemma 4.9. - For $\Lambda_{0}, \Lambda_{1}$ and $M_{2}, M$ above

$$
\begin{gather*}
I_{k} L_{\mathrm{loc}}^{2}\left(M ; \Lambda_{0} \sqcup \Lambda_{1}\right) \stackrel{\beta_{*}}{\longleftrightarrow} I_{k} L_{\nu, \mathrm{loc}}^{2}\left(M_{2}, \vartheta\left(\partial M_{2} \sqcup \beta^{*}\left\{x_{1}=0\right\}\right)\right), \\
I_{k} L_{\mathrm{loc}}^{2}\left(M ; \Lambda_{0}, J_{\Lambda_{0} \cap \Lambda_{1}}^{1} \Lambda_{1}, 2\right) \stackrel{\beta_{*}}{\longleftrightarrow} I_{k} L_{\nu, \mathrm{loc}}^{2}\left(M_{2}, \nabla\left(\partial M_{2}\right)\right), \tag{4.51}
\end{gather*}
$$

where $\beta_{*} \nu=d x$.
Proof. - We only need to check that the left hand side of (4.51) is defined by stability under vector fields - that is clear for the space in the preceding line and the remaining portion of the proof is a straightforward computation. In view of Definition 4.5 we first need to show that

$$
\begin{equation*}
x_{2}^{l} u \in H_{\left(\frac{1}{3} l\right)}(M), \frac{1}{3} l \leq k \Longleftrightarrow u,\left(x_{2}^{3} D_{x_{i}}\right)^{m} u \in L^{2}(M), m \leq k \tag{4.52}
\end{equation*}
$$

The right hand side is equivalent to $\left(x_{2}^{3}\left\langle D_{x}\right\rangle\right)^{m} u \in L^{2}(M)$ and the left hand side to $\left(x_{2}^{3}\left\langle D_{x}\right\rangle^{\frac{1}{3}}\right)^{l} u \in L^{2}(M)$. Since $\left[x_{2},\left\langle D_{x}\right\rangle^{\frac{1}{3}}\right] \in \Psi^{-\frac{2}{3}}(M)$, the left to right implication follows. For the opposite direction we observe that if $u \in C_{0}^{\infty}(M)$ then

$$
\begin{aligned}
\left\|x_{2}\left\langle D_{x}\right\rangle^{\frac{1}{3}} u\right\|_{L^{2}(M)}^{2} & =\int_{M} x_{2}^{2}\left\langle D_{x}\right\rangle^{2 / 3} u \bar{u} d x+\int_{M}\left[\left\langle D_{x}\right\rangle^{1 / 3}, x_{2}^{2}\right]\left\langle D_{x}\right\rangle^{\frac{1}{3}} u \bar{u} d x \\
& \leq\left\|x_{2}^{2}\left\langle D_{x}\right\rangle^{2 / 3} u\right\|_{L^{2}(M)}\|u\|_{L^{2}(M)}+\|u\|_{L^{2}(M)}^{2}
\end{aligned}
$$

and similarly

$$
\left\|x_{2}^{2}\left\langle D_{x}\right\rangle^{\frac{2}{3}} u\right\|_{L^{2}(M)}^{2} \leq\left\|x_{2}^{3}\left\langle D_{x}\right\rangle u\right\|_{L^{2}(M)}\left\|x_{2}\left\langle D_{x}\right\rangle^{\frac{1}{3}} u\right\|_{L^{2}(M)}+\|u\|_{L^{2}(M)}\left\|x_{2}\left\langle D_{x}\right\rangle^{\frac{1}{3}} u\right\|_{L^{2}(M)}
$$

Hence using $2 a b \leq \epsilon a^{2}+\epsilon^{-1} b^{2}$ we conclude that for $u \in C_{0}^{\infty}(M)$

$$
\left\|\left(x_{2}\left\langle D_{x}\right\rangle^{\frac{1}{3}}\right)^{l} u\right\|_{L^{2}(M)} \leq C \sum_{m \leq k}\left\|\left(x_{2}^{3}\left\langle D_{x}\right\rangle u\right)^{m}\right\|_{L^{2}(M)}
$$

if $l \leq 3 k$. A density argument and commutation with the remaining vector fields in the definition of $I_{k} L^{2}\left(M ; \Lambda_{0}, J_{\Lambda_{0} \cap \Lambda_{1}}^{1} \Lambda_{1}, 2\right)$ conclude the proof.

We should remark that the last space on the right hand side is equal to $H_{(k)}^{b}\left(M_{2}\right)$ with the measure $\nu$ and thus the supermarked space on the left is an interpolation space in $k$. For every order $k=2 l$ the space $I_{2 l} L_{\text {loc }}^{2}\left(M ; \Lambda_{1}, \Lambda_{1} \cap \Lambda_{0}, \frac{1}{2}\right)$ is easily seen to be characterized by the condition

$$
\begin{equation*}
\left(x_{1} D_{x_{1}}\right)^{k_{1}}\left(x_{2} D_{x_{2}}\right)^{k_{2}}\left(x_{1} D_{x_{2}}\right)^{k_{3}} D_{x_{2}}^{k_{4}} D_{x^{\prime}}^{\alpha^{\prime}} u \in H_{\left(-\frac{1}{3} k_{4}\right)}(M) \tag{4.53}
\end{equation*}
$$

$\sum_{i<4} k_{i}+\frac{2}{3} k_{4}+\left|\alpha^{\prime}\right| \leq k=2 l$. For $I_{2 l} L_{\text {loc }}^{2}\left(M, \Lambda_{1}^{ \pm}, 2\right), \Lambda_{1}^{ \pm}=\Lambda_{1} \cap\left\{ \pm x_{2} \geq 0\right\}$, we need to add the stability under the operator $\chi\left(\left|D_{x_{1}}\right|\right) \chi\left(\mp\left|D_{x_{1}}\right|^{\frac{1}{3}} x_{2}\right)\left|D_{x_{1}}\right|^{\frac{1}{3}} x_{2}$ with weight $\frac{1}{3}$. The following Lemma constitutes the harder part of the proof of Proposition 4.3 and it will be crucial in chapter 7:

Lemma 4.10. - If

$$
v \in I_{k} L_{\nu, \mathrm{loc}}^{2}\left(M_{2}, \vartheta\left(\beta^{*}\left\{x_{1}=0, \pm x_{2} \geq 0\right\} \sqcup \partial M_{2}\right)\right)
$$

then

$$
\beta_{*} v \in I_{k} L_{\mathrm{loc}}^{2}\left(M ; \Lambda_{1}^{ \pm}, 2\right)+\beta_{*} H_{(k)}^{b}\left(M_{2}\right)
$$

Proof. - We observe first that the two conormal spaces in the Lemma are interpolation spaces in $k$. As the Lagrangian space, $I_{k} L_{\text {loc }}^{2}\left(M ; \Lambda_{1}^{ \pm}, 2\right)$, was defined by complex interpolation for $k$ odd, we only need to prove the lemma for $k=2 l$, $l \in \mathbb{Z}_{+}$. It is also sufficient to consider the + case alone.
We can assume that supp $v \subset \beta^{*}\left\{\left|x_{1}\right|<\epsilon x_{2}^{3}\right\}$, as on the remaining part of the support the function is in $H_{(k)}^{b}\left(M_{2}\right)$, see Fig. 4.1. Thus we can consider supp $v \subset$


$$
M_{2} \xrightarrow{\beta_{1-2}}
$$


$M, \xrightarrow{\beta_{3-1}}$

Figure 4.1. The manifold $M_{2}$ and the support of $v$
$M_{1} \cap \beta_{3-1}^{*}\left\{\left|x_{1}\right|<\epsilon x_{2}^{3}\right\}$ and use coordinates $(X, r), \beta_{3-1}(X, r)=\left(r^{3} X, r\right)$, where we drop the insignificant coordinate $x^{\prime}$. We also have

$$
{ }^{b} T_{\operatorname{supp} v}^{*} M_{2} \simeq\left\{(X, r ; \Xi, \lambda): X, \Xi, \lambda \in \mathbb{R}, r \in \overline{\mathbb{R}}_{+}\right\}
$$

and we can assume that

$$
\begin{equation*}
W F(v) \subset\{(X, r ; \Xi, \lambda):|\lambda|<|\Xi|\} . \tag{4.54}
\end{equation*}
$$

In fact, using a $\Psi_{b}^{0}$-partition of unity, we obtain $v=v_{1}+v_{2}, W F\left(v_{1}\right) \subset\{|\lambda|<|\Xi|\}$ and $W F\left(v_{2}\right) \subset\left\{|\lambda|>\frac{1}{2}|\Xi|\right\}$. Since $\left(r D_{r}\right)^{p} v_{i} \in L_{\nu}^{2}\left(M_{2}\right), p \leq k$ we obtain $v_{2} \in H_{(k)}^{b}\left(M_{2}\right)$, so it can be neglected.
Using the ellipticity of $D_{X}$ in $W F(v)$ we can then write

$$
\begin{equation*}
v=D_{X}^{m} \tilde{v}+v^{\#}, \tilde{v} \in I_{k} H_{(m)}^{b}\left(M_{2}, \mathscr{O}\left(\beta^{*}\left\{x_{1}=0, x_{2}>0\right\} \sqcup \partial M_{1}\right)\right), \quad v^{\#} \in H_{(k)}^{b}\left(M_{2}\right) \tag{4.55}
\end{equation*}
$$

Replacing $v$ by $v-v^{\#}$, we now claim that $\beta_{*} v \in I_{k} L_{\text {loc }}^{2}\left(M ; \Lambda_{1}^{+}, 2\right)$. Since all the other operators in (4.53) lift we only need to investigate $D_{x_{2}}^{p}$ and

$$
\left.\chi\left(\left|D_{x_{1}}\right|\right) \chi\left(\mp \left\lvert\, D_{x_{1}}{ }^{\frac{1}{3}} x_{2}\right.\right) \right\rvert\, D_{x_{1}}{ }^{\frac{1}{3}} x_{2}
$$

Thus, we start with $p=3 m \leq 2 l=k$ and

$$
\begin{align*}
& D_{x_{2}}^{3_{m}} \beta_{*} v=\beta_{*}\left(\left\{\frac{1}{r}\left(r D_{r}-3 X D_{X}\right)\right\}^{3 m} D_{X}^{m} \tilde{v}\right)= \\
& =D_{x_{1}}^{m} \beta_{*}\left(r^{3 m}\left\{\frac{1}{r}\left(r D_{r}-3 X D_{x}\right)\right\}^{3 m} \tilde{v}\right) \\
& \quad+\beta_{*}\left(\left[r^{3 m}\left\{\frac{1}{r}\left(r D_{r}-3 X D_{X}\right)\right\}^{3 m}, r^{-3 m} D_{X}^{m}\right] \tilde{v}\right) \tag{4.56}
\end{align*}
$$

where we used $D_{x_{1}} \beta_{*}=\beta_{*} r^{-3} D_{X}$. Since $\left[r^{-1}\left(r D_{r}-3 X D_{X}\right), D_{X}\right]=r^{-1} 3 i D_{X}$, the expansion of the commutator gives

$$
D_{x_{2}}^{3 m} \beta_{*} v=D_{x_{1}}^{m} \beta_{*} v_{1}, v_{1} \in I_{2 k-2 m} L_{\nu, \text { loc }}^{2}\left(M_{2}, \vartheta\left(\beta^{*}\left\{x_{1}=0, x_{2}>0\right\} \sqcup \partial M_{2}\right)\right),
$$

so that $D_{x_{2}}^{3 m} \beta_{*} v \in H_{(-m)}(M)$ if $2 m \leq 2 l=k$, that is $\frac{2}{3} p \leq k, p=3 m$. Thus the weight of $D_{x_{2}}^{p}$ is in agreement with the weights in $I_{k} L_{\text {loc }}^{2}\left(M, \Lambda_{1}^{+}, 2\right)$. For $p$ not divisible by 3 we use

$$
\begin{align*}
\left\|D_{x_{2}} u\right\|_{H_{\left(-\frac{1}{3}\right)(M)}}^{2} & \leq \frac{1}{3}\left\|D_{x_{2}}^{3} u\right\|_{H_{(-1)}(M)}^{2}+\frac{2}{3}\|u\|_{L^{2}(M)}^{2} \\
\left\|D_{x_{2}}^{2} u\right\|_{H_{\left(-\frac{2}{3}\right)}(M)}^{2} & \leq \frac{2}{3}\left\|D_{x_{2}}^{3} u\right\|_{H_{(-1)}(M)}^{2}+\frac{1}{3}\|u\|_{L^{2}(M)}^{2}, u \in C_{0}^{\infty}(M) \tag{4.57}
\end{align*}
$$

In (4.53), $D_{x_{2}}$ is the only operator with a fractional weight $\left(\frac{2}{3}\right)$, so that if $D_{x_{2}}^{2}$ occurs we can use the stability under $D_{x_{2}}^{3}$. If $V D_{x_{2}} A$ occurs, where $A$ is a product of operators and $V$ is an operator different from $D_{x_{2}}$, then we can use (4.57), integration by parts and stability under $V^{2} A, D_{x_{2}}^{3} A,\left(D_{x_{2}}^{3} A \in H_{(-1)}(M)\right)$. This proves (4.53) for $\beta^{*} v$ and we still need to consider the remaining multiplier. That however is easy, since for $v$, after all the reductions above $\left(\bmod H_{(k)}^{b}\left(M_{2}\right)\right)$,

$$
\left.\chi\left(\left|D_{x_{1}}\right|\right) \chi\left(-\left|D_{x_{1}}\right|^{\frac{1}{3}} x_{2}\right) \right\rvert\, D_{x_{1}}{ }^{\frac{1}{3}} x_{2} \beta_{*} v=0 .
$$

In fact, let us consider $H \in L^{\infty}\left(M_{2}\right)$ defined by

$$
\left.H\right|_{M_{2} \backslash \partial M_{2}}=\left(\left.\beta\right|_{M_{2} \backslash \partial M_{2}}\right)^{*}\left(\left(x_{2}\right)_{-}^{0}\right) .
$$

Then $\beta_{*} H=\left(x_{2}\right)_{-}^{0}$ and $\left(x_{2}\right)_{-}^{0} \beta_{*} v=\beta_{*} H v=0$, if as assumed, $\operatorname{supp} v \subset \beta^{*}\left\{\left|x_{1}\right|<\right.$ $\left.\epsilon x_{2}^{3}\right\}$. We then observe that

$$
\left.\chi\left(\left|D_{x_{1}}\right|\right) \chi\left(-\left|D_{x_{1}}\right|^{\frac{1}{3}} x_{2}\right)\left|D_{x_{1}}{ }^{\frac{1}{3}} x_{2}\left(x_{2}\right)_{-}^{0}=\chi\left(\left|D_{x_{1}}\right|\right) \chi\left(-\left|D_{x_{1}}\right|^{\frac{1}{3}} x_{2}\right)\right| D_{x_{1}}\right|^{\frac{1}{3}} x_{2}
$$

concluding the proof.
Remark 4.11. - A converse of Lemma 4.2 is also true:

$$
\beta^{*} I_{k} L_{\mathrm{loc}}^{2}\left(M ; \Lambda_{1}^{ \pm}, 2\right) \subset I_{k} L_{\nu, \mathrm{loc}}^{2}\left(M_{2}, \mathscr{O}\left(\beta^{*}\left\{x_{1}=0, \pm x_{2} \geq 0\right\} \sqcup \partial M_{2}\right)\right),
$$

and although we do not need it, this fact is implicitly present in the Dirichlet estimates. Its failure (see Remark 7.8) in one case considered in chapter 7 explains to some extent the complications of the space $J_{k} L^{2}(X)$.

Proof. - Proof of Proposition 4.3 It is easy to see that the right hand sides in (4.47) and (4.48) are contained in $I_{k} L_{\text {loc }}^{2}(M ; \Lambda \sqcup \widetilde{\Lambda})$. Since $I_{k} L^{2}\left(M ; \mathfrak{\varrho}^{ \pm}\right) \subset I_{k} L^{2}\left(M, \mathfrak{£}_{\frac{1}{2}}\right)$, it remains to verify (4.48) and that easily follows from Lemmas 4.1 and 4.2 as

$$
\begin{aligned}
& I_{k} L_{\nu, \mathrm{loc}}^{2}\left(M_{2}, \circlearrowleft\left(\partial M_{2} \sqcup \beta^{*}\left\{x_{1}=0\right\}\right)\right)= \\
& H_{(k)}^{b}\left(M_{2}\right)+\sum_{ \pm} I_{k} L_{\nu, \mathrm{loc}}^{2}\left(M_{2}, \circlearrowleft\left(\partial M_{2} \sqcup \beta^{*}\left\{x_{1}=0, \pm x_{2}>0\right\}\right)\right),
\end{aligned}
$$

which concludes the proof.

Remark 4.12. - The same proof can be used for $t=1$, with the $3-1$ and $1-2$ blow-up replaced by a 2-1 and 1-1 blow-ups of the same submanifolds (with the obvious modification in the latter case). The last ( $1-1$ ) blow-up is not strictly speaking necessary but it is useful in characterizing $\beta^{*} I_{k}\left(M ; \Lambda_{0} \sqcup \Lambda_{1}\right)$ as a b-Sobolev space.
4.3. Let $M$ be a $C^{\infty}$ manifold with a $C^{\infty}$ boundary $\partial M$. The $b$-Lagrangian distributions in Sect. 7 of [34] were defined for $C^{\infty}$ homogeneous Lagrangian submanifolds $\Lambda \subset T^{*} M \backslash 0$ such that $\jmath \Lambda \subset{ }^{b} T^{*} M \backslash 0$ is smooth,

$$
\jmath: T^{*} M \backslash 0 \longrightarrow{ }^{b} T^{*} M \backslash 0 .
$$

They were defined using (4.2) and (4.3) with totally characteristic operators, $\Psi_{b}^{1}(M)$ employed in place of $\Psi^{1}$. We need to generalize this notion to allow cases when $\jmath \Lambda \subset{ }^{b} T^{*} M \backslash 0$ is singular in a controlled way. A $C^{\infty}$ submanifold of ${ }^{b} T^{*} M \backslash 0,{ }^{b} \Lambda$, is called a homogeneous (b-)Lagrangian if and only if $\operatorname{dim}^{b} \Lambda=\operatorname{dim} M,{ }^{b} \Lambda \cap^{b} T_{\partial M}^{*} M$ is a homogeneous $C^{\infty}$ Lagrangian submanifold of $T^{*} \partial M$,

$$
\left.{ }^{b} \omega\right|_{T_{m}{ }^{b} \Lambda}=0 \quad \forall m \in{ }^{b} \Lambda \cap^{b} T_{M^{\circ}}^{*} M \backslash 0,
$$

and ${ }^{b} \Lambda$ is homogeneous with respect to the natural $\mathbb{R}_{+}$-action on ${ }^{b} T^{*} M \backslash 0$. The last two conditions simply state that ${ }^{b} \Lambda \cap{ }^{b} T_{M^{\circ}}^{*} M \backslash 0 \subset T^{*} M^{\circ} \backslash 0$ is a homogeneous Lagrangian submanifold of $T^{*} M^{\circ} \backslash 0$. Locally ${ }^{b} \Lambda$ is given by the zeros of $n=\operatorname{dim} M$ functions in $C^{\infty}\left({ }^{b} T^{*} M \backslash 0\right)$ :

$$
\begin{gathered}
{ }^{b} \Lambda \cap \Gamma=\left\{m \in \Gamma: f_{1}(m)=\cdots=f_{n}(m)=0\right\} \\
f_{i} \in C^{\infty}\left({ }^{b} T^{*} M \backslash 0\right), \quad{ }^{b}\left\{f_{i}, f_{j}\right\}=0,
\end{gathered}
$$

where $\Gamma \subset{ }^{b} T^{*} M \backslash 0$ is a conic neighbourhood of $\bar{m} \in{ }^{b} \Lambda$. To pass to the global situation we consider the ideal

$$
\mathscr{y}_{b_{\Lambda}}=\left\{f \in C^{\infty}\left({ }^{b} T^{*} M \backslash 0\right):\left.f\right|_{b_{\Lambda=0}}\right\}
$$

which is locally finitely generated. Since ${ }^{b} \Lambda$ is homogeneous we introduce for every $k \in \mathbb{Z}$

$$
{ }^{b_{9}} \pi_{k}\left({ }^{b} \Lambda\right)=\mathscr{G}_{\Lambda^{b}} \cap S_{\mathrm{hg}}^{k}\left({ }^{b} T^{*} M \backslash 0\right)
$$

and observe that

$$
{ }^{b} \Lambda=\left\{m: f(m)=0, \forall f \in \in^{b} \Re_{k}\left({ }^{b} \Lambda\right)\right\}
$$

Example 4.1. If $\Lambda_{0} \subset T^{*} \partial M \backslash 0$ is a homogeneous Lagrangian submanifold and $M \simeq \partial M \times[0,1)$, then $\Lambda=\left\{(x, y ; \xi, \eta):(y, \eta) \in \Lambda_{0}, \xi=0\right\} \subset T^{*} M \backslash 0$ is a Lagrangian submanifold and $\jmath \Lambda \subset{ }^{b} T^{*} M \backslash 0$ is a $C^{\infty} b$-Lagrangian submanifold.
Example 4.2 Let us define $\tilde{\Lambda} \subset T^{*} \mathbb{R}^{2} \backslash 0$ by $\tilde{\Lambda}=N^{*}\left\{x+y^{2}=0\right\}$. If $M=\overline{\mathbb{R}}_{+}^{2}=$ $\{x \geq 0\}$ and $\Lambda=\left.\tilde{\Lambda}\right|_{T^{*} M \backslash 0}$ then ${ }^{b} T^{*} M \backslash 0 \supset \jmath \Lambda=\left\{(x, y ; \lambda, \eta): x+y^{2}=\eta x-2 y \lambda=0\right.$, $(\eta, \lambda) \neq(0,0)\}$ is not smooth.

We wish to consider Lagrangian submanifolds similar to those in Example 4.2 in the sense of having defining functions which are polynomials in $\xi$ at the boundary in a way which is invariant under $b$-canonical transformations - see [26], part III.

We use $x$ as the defining function of the boundary and $\lambda={ }^{b} \sigma_{1}\left(x D_{x}\right)$ as its $b$-dual variable in ${ }^{b} T^{*} M \backslash 0$. Then we consider

$$
\begin{equation*}
S_{\mathrm{hg}}^{m, k}\left({ }^{b} T^{*} M \backslash 0\right)=\left\{\sum_{0 \leq p \leq k} \frac{\lambda^{p}}{x^{p}} a_{p}: a_{p} \in S_{\mathrm{hg}}^{m-p}\left({ }^{b} T^{*} M \backslash 0\right)\right\}, \tag{4.58}
\end{equation*}
$$

where $S_{\mathrm{hg}}^{l}\left({ }^{b} T^{*} M \backslash 0\right)$ are homogeneous symbols $a \in C^{\infty}\left({ }^{b} T^{*} M \backslash 0\right)$. These are the symbols satisfying $a\left(T_{r} m\right)=r^{l} a(m)$, where $T_{r}$ is the $\mathbb{R}_{+}$action generated by $\rho_{\partial M}+\lambda \partial_{\lambda}, \rho_{\partial M}$ the radial vector field on $T^{*} \partial M \backslash 0$.

The right-hand side of (4.58) involves a choice of coordinates but the space $S_{\mathrm{hg}}^{m, k}$ is nevertheless invariantly defined. In fact, for $k=0, S_{\mathrm{hg}}^{m, k}=S_{\mathrm{hg}}^{m}$ and $S_{\mathrm{hg}}^{m, 1}$ can be characterized by demanding that $\tilde{x} \cdot S_{\mathrm{hg}}^{m, 1} \subset S_{\mathrm{hg}}^{m}$ and that $\left.\tilde{x} \cdot S_{\mathrm{hg}}^{m-1}\right|_{{ }_{b} T_{\partial M}^{*} M}$ vanishes at $\tilde{\lambda}={ }^{b} \sigma_{1}\left(\tilde{x} D_{\tilde{x}}\right)=0$, for any defining function of $\partial M$. Since $\tilde{\lambda}$ is well defined on ${ }^{b} T_{\partial M}^{*} M$ this gives an invariant definition. We then see that

$$
S_{\mathrm{hg}}^{m, k}=\left\{\sum_{0 \leq p \leq k} b_{l_{1}} \ldots b_{l_{p}}: b_{l_{i}} \in S_{\mathrm{hg}}^{m, 1}\right\}
$$

is also invariantly defined. We should also remark that

$$
\begin{equation*}
\left.S_{\mathrm{hg}}^{m, k}\left({ }^{b} T^{*} M \backslash 0\right)\right|_{T^{*} M^{\circ} \backslash 0}=\left\{\left.\sum_{0 \leq o \leq k} \xi^{p} \iota^{*} a_{p}\right|_{T^{*} M^{\circ} \backslash 0}: a_{p} \in S_{\mathrm{hg}}^{m-p}\left({ }^{b} T^{*} M \backslash 0\right)\right\} \tag{4.59}
\end{equation*}
$$

where $\xi=\sigma_{1}\left(D_{x}\right)$. This characterizes $S_{\mathrm{hg}}^{m, k}$ as well. Clearly, the homogeneous symbols ( $S_{\mathrm{hg}}$ ) could be replaced by the polyhomogeneous ones ( $S_{\mathrm{phg}}$ ), or by the usual symbols.

Definition 4.13. - A homogeneous $C^{\infty}$ Lagrangian submanifold $\Lambda \subset T^{*} M \backslash 0$ is called b-polynomially defined if and only if there exists an ideal $\mathcal{G} \subset \bigcup_{m, k} S_{\mathrm{phg}}^{m, k}\left({ }^{b} T^{*} M \backslash 0\right)$ such that

$$
\Lambda=\left\{m \in T^{*} M \backslash 0: \jmath^{*} f(m)=0, \quad \forall f \in \mathscr{g}\right\}
$$

Putting

$$
\begin{equation*}
{ }^{b_{\mathfrak{G}}}{ }_{p}^{\#}(\Lambda)=\mathscr{g} \cap \bigcup_{k \geq 0} S_{\mathrm{phg}}^{p, k}\left({ }^{b} T^{*} M \backslash 0\right) \tag{4.60}
\end{equation*}
$$

and using elliptic elements of $S_{\mathrm{phg}}^{0}\left({ }^{b} T^{*} M \backslash 0\right)$ we see that a $b$-polynomially defined Lagrangian $\Lambda$ is given as

$$
\begin{equation*}
\Lambda=\left\{m \in T^{*} M \backslash 0: \jmath^{*} a(m)=0, \forall a \in^{b} \mathfrak{R}_{p}^{\#}(\Lambda)\right\} \tag{4.61}
\end{equation*}
$$

We can define ${ }^{b} \mathfrak{M}_{p}^{\#}(\Lambda)$ for an arbitrary $\Lambda \subset T^{*} M \backslash 0$ as the set of $a \in$ $\bigcup_{k \geq 0} S_{\mathrm{phg}}^{p, k}\left({ }^{b} T^{*} M \backslash 0\right)$ such that $\left.j^{*} a\right|_{\Lambda}=0$. Then for $b$-polynomially defined Lagrangians
(4.61) holds. One easily observes that if $\jmath \Lambda={ }^{b} \Lambda$ is a smooth $b$-Lagrangian, then ${ }^{b} \mathfrak{M}_{m}^{\#}(\Lambda) \supset^{b} \mathscr{R}_{1}(\Lambda)$ with strict inclusion in general (see for instance Example 4.1 above; equality occurs for $\left.\Lambda=N^{*} \partial M\right)$.

The property of being $b$-polynomially defined is local near each fiber of $T_{\partial M}^{*} M$. Thus if $m \in T_{\partial M}^{*} M$, we say a $C^{\infty}$ Lagrangian $\Lambda$ is b-polynomially defined near $m$ if for an open cone $\Gamma \subset T^{*} M \backslash 0, m \in \Gamma$

$$
\begin{equation*}
\Lambda \cap \Gamma=\left\{m \in \Gamma: \jmath^{*} f(m), \forall f \in \mathfrak{L}_{\Gamma}\right\} \tag{4.62}
\end{equation*}
$$

where $\mathscr{I}_{\Gamma}$ is an ideal in $\bigcup_{p, k} S_{\mathrm{phg}}^{p, k}\left({ }^{b} T^{*} M \backslash 0\right)$. If $\Gamma$ is everywhere locally $b$-polynomially defined, a partition of unity argument shows that $\Lambda$ is $b$-polynomially defined. We could rephrase the local definition by saying that $\Lambda$ is defined by ${ }^{b} \mathfrak{M}_{p}^{\#}(\Lambda)$ near $m$ (or in $\Gamma$ ).
Example 4.3. We give an example of a Lagrangian submanifold which is not $b$ polynomially defined and to do that we start with the non-homogeneous case. Thus we define $\Lambda \subset N \subset T^{*} \overline{\mathbb{R}}_{+}^{2}$, for a small neighbourhood of $(0,0)$, $N$, as the set of zeros of

$$
\begin{aligned}
& f_{1}(x, y ; \xi, \eta)=f(\xi)+y(\xi+1) \\
& f_{2}(x, y ; \xi, \eta)=x+\eta(\xi+1)^{-1}\left(y+f^{\prime}(\xi)\right)
\end{aligned}
$$

where $f \in C^{\infty}(\mathbb{R}), f^{(k)}(0)=0$ for all $k$ and $f(t) \neq 0$ if $t \neq 0$. We easily check that $d f_{1}(0 ; 0)$ and $d f_{2}(0 ; 0)$ are linearly independent and that $\left\{f_{1}, f_{2}\right\} \equiv 0$. Since $\partial_{\xi}^{k} f_{1}(0 ; 0)=0$ for all $k, f_{1}(0, y ; \xi, \eta)$ cannot be written as a polynomial in $\xi$. By introducing an additional variable $z$ with the dual $\zeta$, a homogeneous example is obtained by taking $\bar{f}_{i}(x, y, z ; \xi, \eta, \zeta)=f_{i}(x, y ; \xi / \zeta, \eta / \zeta), i=1,2$ and choosing a homogeneous function $\bar{f}_{3}$, so that $\left\{\bar{f}_{i}, \bar{f}_{3}\right\}=0$, as we may by Darboux's theorem. The vanishing of $\bar{f}_{i}$ 's defines a smooth homogeneous Lagrangian submanifold in a conic neighbourhood of $(0 ;(0,0,1)) \in T^{*} \mathbb{R}^{3} \backslash 0$.

If $\Gamma^{0}=\Gamma \cap T^{*} M^{\circ} \backslash 0$, where $\Gamma$ is the open cone above, we can consider a canonical transformation

$$
\begin{equation*}
\chi: \Gamma^{\circ} \longrightarrow T^{*} M^{\circ} \backslash 0 \tag{4.63}
\end{equation*}
$$

which we assume extends to a b-canonical transformation [25]

$$
\begin{equation*}
{ }^{b} \chi:{ }^{b} \Gamma \longrightarrow{ }^{b} T^{*} M \backslash 0, \quad \jmath(\Gamma) \subset{ }^{b} \Gamma \tag{4.64}
\end{equation*}
$$

such that

$$
\begin{equation*}
{ }^{b} \chi(\jmath(m))=\jmath\left(m_{1}\right), \quad m_{1} \in T_{\partial M}^{*} M \backslash 0 \tag{4.65}
\end{equation*}
$$

The basic invariance property is now given in

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Lemma 4.14. - Let us assume that $\Lambda, \Lambda_{1} \subset T^{*} M \backslash 0$ are homogeneous $C^{\infty}$ Lagrangian submanifolds and $\chi$ with the properties (4.63), (4.64) and (4.65) satisfies also

$$
\begin{equation*}
\chi\left(\Lambda \cap \Gamma^{0}\right) \subset \Lambda_{1}, \quad N^{*} \partial M \cap \Lambda=\emptyset \tag{4.66}
\end{equation*}
$$

If $\Lambda_{1}$ is locally b-polynomially defined near $m_{1}$, then $\Lambda$ is locally b-polynomially defined near $m$.

Proof. - The second condition in (4.66) shows that $\Lambda=\operatorname{cl}\left(\Lambda \cap T^{*} M^{0} \backslash 0\right)$ which in view of the first condition holds also for $\Lambda_{1}$. Thus $\Lambda_{1}$ is defined by $\left.{ }^{b} \Omega_{1}^{\#}\left(\Lambda_{1}\right)\right|_{T^{*} M^{\circ} \backslash 0}$ near $m_{1}$. Since $\left.{ }^{b} \Re_{1}^{\#}\left(\Lambda_{1}\right)\right|_{T^{*} M^{\circ} \backslash 0} \subset S^{1}\left(T^{*} M^{\circ} \backslash 0\right), \Lambda$ is defined by $\chi^{*}\left(\left.{ }^{b} 9 \Re_{1}^{\#}\left(\Lambda_{1}\right)\right|_{T^{*} M^{\circ} \backslash 0}\right)$ near $m$. Recalling that ${ }^{b} \chi^{*} x=a x$ and ${ }^{b} \chi^{*} \lambda=b \lambda+c x, a, b \neq 0$, and that ${ }^{b} \chi^{*}$ preserves $S_{\mathrm{phg}}^{m}\left({ }^{b} T^{*} M \backslash\right)$, we easily see that

$$
{ }^{b} \chi^{*}: S_{\mathrm{phg}}^{m, k}\left({ }^{b} T^{*} M \backslash 0\right) \longrightarrow S_{\mathrm{phg}}^{m, k}\left({ }^{b} \Gamma\right)
$$

Thus $\chi^{*}\left(\left.{ }^{b} 9 \mathbb{R}_{1}^{\#}\left(\Lambda_{1}\right)\right|_{T^{*} M^{0} \backslash 0}\right)$ extends to $\bigcup_{k \geq 0} S_{\mathrm{phg}}^{m, k}\left({ }^{b} \Gamma\right)$, and this extension defines $\Lambda$ near $m$.

We next consider marked Lagrangians. Let $K \subset \Lambda$ be a homogeneous submanifold of codimension 1 satisfying the condition
$K$ is tangent to $T_{\partial X}^{*} X$ to a fixed finite order.
Note that we do not demand that $K=\{m \in \Lambda: a(m)=0\}, a \in S_{\mathrm{phg}}^{m, k}\left({ }^{b} T^{*} M \backslash 0\right)$ even though this is the case in all applications. For $\Lambda$ and $K$ satisfying (4.67) we define,

$$
{ }^{b} \mathfrak{N}_{1}^{\#}(\Lambda, K)=\left\{a \in^{b} \mathfrak{M}_{1}^{\#}(\Lambda): H_{\left(\left.\jmath^{*} a\right|_{T^{*} M^{\circ} \backslash 0}\right)} \text { is tangent to } K \cap T^{*} M^{\circ} \backslash 0\right\} .
$$

The $b$-canonical transformations preserve $b$-Hamilton vector fields which satisfy

$$
\left.{ }^{b} H_{f}\right|_{\left({ }^{b} T^{*} M \backslash 0\right)^{\circ}}=H_{\left.J^{*} f\right|_{T^{*} M^{\circ} \backslash 0}} .
$$

Thus, because of (4.64), the definition above is invariant and we obtain the following analogue of Lemma 4.3:

Lemma 4.15. - Let $\chi, \Lambda, \Lambda, \Lambda_{1}$ be as in Lemma 4.3 and assume that homogeneous hypersurfaces $K \subset \Lambda, K_{1} \subset \Lambda$ satisfy (4.67) and

$$
\begin{equation*}
\chi(K \cap \Gamma) \subset K_{1} \tag{4.68}
\end{equation*}
$$

Then

$$
f \in{ }^{b} \mathfrak{M}_{1}^{\#}\left(\Lambda_{1}, K_{1}\right), \operatorname{supp} f \subset{ }^{b} \chi\left({ }^{b} \Gamma\right) \Longrightarrow{ }^{b} \chi_{1}^{*} f \in \mathfrak{M}_{1}^{\#}(\Lambda, K)
$$

To quantize these geometric notions we recall the definition of a class of operators from [25], Sect.III.5:

$$
\Psi_{b}^{m, k}(M)=\left\{\sum_{\text {finite }} P_{i} A_{i}+B_{i} Q_{i}: \quad A_{i}, B_{i} \in \Psi_{b}^{m-m_{i}}(M), P_{i}, Q_{i} \in \operatorname{Diff}^{m_{i}}(M), m_{i} \leq k\right\}
$$

By restricting to the interior $M^{\circ}$, and using (4.59) one easily obtains a surjective symbol map

$$
{ }^{b} \sigma_{m, k}: \Psi_{b}^{m, k}(M) \rightarrow S^{m, k}\left({ }^{b} T^{*} M \backslash 0\right), \quad{ }^{b} \sigma_{m, k}(P A)=\sigma_{\operatorname{ord}(P)}(P) \cdot{ }^{b} \sigma_{m-\operatorname{ord}(P)}(A)
$$

if $P$ is a homogeneous differential operator of $\operatorname{ord}(P) \leq k$. These leads us to define

$$
\begin{equation*}
\Psi_{b}^{m, k}(M, \Lambda)=\left\{C \in \Psi_{b}^{m, k}(M):{ }^{b} \sigma_{m, k}(C) \in{ }^{b} 9 \mathbb{m}_{m}^{\#}(\Lambda)\right\} \tag{4.69}
\end{equation*}
$$

and similarly

$$
\Psi_{b}^{m, k}(M ; \Lambda, K)=\left\{C \in \Psi_{b}^{m, k}(M):{ }^{b} \sigma_{m, k}(C) \in^{b} \mathfrak{M}_{m}^{\#}(\Lambda, K)\right\}
$$

We can also consider a disjoint union of two cleanly intersecting Lagrangians $\Lambda_{1}, \Lambda_{2}$, $\Lambda_{1} \sqcup \Lambda_{2}$ and define ${ }^{b} 9 \Re_{m}^{\#}\left(\Lambda_{1} \sqcup \Lambda_{2}\right)$ and $\Psi_{b}^{m, k}\left(M, \Lambda_{1} \sqcup \Lambda_{2}\right)$ with the analogue of Lemma 4.4 easily available. We note that

$$
\begin{align*}
& b_{\mathfrak{G}}^{\mathbb{R}_{m}^{\#}}\left(\Lambda_{1} \sqcup \Lambda_{2}\right) \subset{ }^{b_{\mathfrak{~}}} \mathbb{R}_{m}^{\#}\left(\Lambda_{i}\right), \quad i=1,2 \\
& \Psi_{b}^{m, k}\left(\Lambda_{1} \sqcup \Lambda_{2}\right) \subset \Psi_{b}^{m, k}\left(\Lambda_{i}\right) . \tag{4.70}
\end{align*}
$$

Let us consider the following Lagrangian varieties

$$
\begin{equation*}
\mathfrak{L}=\Lambda, \quad\{\Lambda, K\} \quad \text { or } \quad \Lambda_{1} \sqcup \Lambda_{2} \tag{4.71}
\end{equation*}
$$

with $\Lambda, K, \Lambda_{i}$ above. We say that $\mathfrak{L}$ is $b$-polynomially defined if the Lagrangians in $\mathfrak{£}$ are $b$-polynomially defined and $K$ satisfies (4.67). We then have

Definition 4.16. - For a Lagrangian variety $\mathfrak{L}, l \in \mathbb{N}_{0}$ and $s \in \mathbb{R}$, we define the space of distributions

$$
\begin{align*}
& I_{l}^{b} \bar{H}_{(s)}(M ; \mathfrak{\mathscr { L }})=\left\{u \in \bar{H}_{(s)}(M): C_{1} \ldots C_{l^{\prime}} u \in \bar{H}_{(s)}(M)\right. \\
& \left.\quad \text { if } C_{i} \in \Psi_{b}^{1, k}(M, \mathfrak{£}) \text { for some } k, \text { and any } l^{\prime} \leq l\right\} \tag{4.72}
\end{align*}
$$

To study the invariance we recall from [25], Sect.III.5, that if $F$ is an elliptic $b$-Fourier Integral Operator of order 0 and $G$ its parametrix then

$$
\begin{equation*}
F \Psi_{b}^{m, k}(M) G \subset \Psi_{b}^{m, k}(M) \tag{4.73}
\end{equation*}
$$

If $F$ is associated to a $b$-canonical transformation satisfying (4.63) and (4.64) we rewrite (4.66) and (4.68) as

$$
\begin{equation*}
\chi(\mathscr{L} \cap \Gamma) \subset \mathscr{L}_{1} \tag{4.74}
\end{equation*}
$$

with an obvious modification for $\mathfrak{L}=\Lambda_{1} \sqcup \Lambda_{2}$. Let $C \in \Psi_{b}^{m, k}(M)$ and $C_{1}=F \circ C \circ G$. By (4.73), $C_{1} \in \Psi_{b}^{m, k}(M)$ and we want to compute its symbol. Let $F^{0}, G^{0}, C^{0}$ and $C_{1}^{0}$ be restriction of the operators above to $M^{\circ}$. Then $C^{0}$ and $C_{1}^{0}$ are pseudodifferential operators of order $m$ on compact subsets of $M^{\circ}$. Egorov's Theorem (see Theorem 25.3.5 of [14]) shows that $\chi^{*} \sigma_{m}\left(F^{0} \circ C^{0} \circ G^{0}\right)=\sigma_{m}\left(C^{0}\right)$ on the subsets of $\Gamma$ which have compact projections to $M^{\circ}$. It follows from the definition of ${ }^{b} \sigma_{m, k}$ that

$$
{ }^{b} \sigma_{m, k}\left(C_{1}\right)={ }^{b} \chi^{* b} \sigma_{m, k}\left(F \circ C_{1} \circ G\right)={ }^{b} \sigma_{m, k}(C) \quad \text { in }{ }^{b} \Gamma .
$$

Combining this discussion with Lemmas 4.3 and 4.4 we obtain
Proposition 4.17. - If $\chi$ satisfies (4.63), (4.64) and (4.74) and $F$ is an elliptic $b$ Fourier Integral Operator of order 0 associated to ${ }^{b} \chi$, then for $£$ locally b-polynomially defined near $\Gamma$

$$
u \in I_{l}^{b} \bar{H}_{(s)}(M, \mathfrak{L}) \quad \text { and } W F_{b}(u) \subset{ }^{b} \Gamma \Longrightarrow F u \in I_{l}^{b} \bar{H}_{(s)}(M, \mathfrak{£})
$$

The main application will be to Lagrangians which are actually polynomially defined. That property does not carry the necessary invariance but it is very useful for computations (see chapter 6).

Proposition 4.18. - Suppose that $\mathfrak{\&}$ is b-polynomially defined and that ${ }^{b} ๑_{1}^{\#}(\mathfrak{L})$ is spanned over $\Psi_{b}^{0}(M)$ by $\left\{A_{j} Q_{j}\right\}_{j \in J}, Q_{j} \in \operatorname{Diff}^{m_{j}}(M), A_{j} \in \Psi_{b}^{1-m_{j}}(M)$, $A_{j}$ elliptic (in the totally characteristic sense). Then

$$
\begin{equation*}
u \in I_{k} \bar{H}_{(s)}(M, \mathscr{\varrho}) \Longrightarrow Q_{j_{1}} \ldots Q_{j_{l}} u \in \bar{H}_{\left(s+l-\sum_{p \leq l} m_{\left.j_{p}\right)}\right)}(M), l \leq k, \quad j_{p} \in J . \tag{4.75}
\end{equation*}
$$

Proof. - Definition 4.16 implies that $A_{j_{1}} Q_{j_{1}} \ldots A_{j_{l}} Q_{j_{l}} u \in \bar{H}_{(s)}(M)$. We recall the commutation relation:

$$
\begin{align*}
& Q \in \Psi_{b}^{m, k}(M), A \in \Psi_{b}^{l}(M) \Longrightarrow A Q=Q A_{1}+Q^{\prime} \\
& A_{1} \in \Psi_{b}^{l}(m),{ }^{b} \sigma_{l}(A)={ }^{b} \sigma_{l}\left(A_{1}\right), Q^{\prime} \in \Psi_{b}^{m+l-1, k} \tag{4.76}
\end{align*}
$$

which follows from iterating

$$
A V=V\left(x A x^{-1}\right)+[V x, A] x^{-1}, V \in \operatorname{Diff}^{1}(M), \quad[V x, A] x^{-1} \in \Psi_{b}^{l-1}(M)
$$

We also observe that

$$
\Psi_{b}^{m, k}(M) \cdot I_{q} \bar{H}_{(r)}(M, \partial M) \subset I_{q-m+k} \bar{H}_{(r-k)}(M, \partial M), q, q-m+k \geq 0 .
$$

Thus by successive applications of (4.76), we obtain that for an elliptic operator $A \in \Psi_{b}^{L}(M), L=l-\sum_{p \leq l} m_{j_{p}}, A Q_{j_{1}} \ldots Q_{j_{l}} u \in I_{-L} \bar{H}_{(s+L)}(M, \partial M)$. Let $B$ be the parametrix of $A$ in $\Psi_{b}^{-L}(M)$. Applying it to the previous expression we obtain (4.75).

## 5. REFINED ESTIMATES IN THE PAST

In the space of distributions defined in chapter 3 only the forward cone and part of the reflected front (the cusp) appear. Since the only data in the past is the incident front this is natural in the case of forward propagation. However, in principle some singularities could appear in the past on the boundary of the resolved space $\widetilde{X}_{4}$ and the purpose of this section is to present some refined conormal estimates which show the absence of such a behaviour. The essential component here is the use of weighted $L^{2}$-estimates 'in the past' similar to those in [32].

We will use the coordinates introduced in Propositions 3.1 and 3.2 and define the non-homogeneous past, $K$, using the model operator $P_{0}$. To start, we define $q_{0} \in C^{\infty}\left(\widetilde{X}_{1}\right)$, independent of $r$ and $y$ as

$$
q_{0}(\omega)=r^{-4}\left(\beta^{*} q\right)(r, \omega), \quad q(x)=x_{1}^{4}+8 x_{1} x_{2}-6 x_{1}^{2} x_{2}-3 x_{2}^{2}
$$

Thus $\beta^{*} Q=\left\{(r, \omega, y): q_{0}(\omega)=0\right\}$ is the model cone (see Proposition 3.1). Throughout this section we will write $\beta \equiv \beta_{1}$. We then consider

$$
\begin{equation*}
\widetilde{X}_{1} \backslash\left\{(r, \omega, y): q_{0}(\omega)=-\epsilon\right\}, \quad \epsilon>0 \tag{5.1}
\end{equation*}
$$

which for small $\epsilon$ has three components.
Definition 5.1. - The past in $\widetilde{X}_{1}, K_{1}$, is defined as the component of (5.1) which contains $\beta^{*} Q_{-}$, where $Q_{-}$is the retarded model cone over $\Gamma$. We then define

$$
\begin{equation*}
K=\beta_{*} K_{1} \tag{5.2}
\end{equation*}
$$

The intersection of $K_{1}$ and $\partial \widetilde{X}_{1}$ is shown in Fig. 5.1. It is actually convenient to consider also a 1-2-3 homogeneous change of variables (cf. Proposition 3.3) which allows us to write

$$
\begin{equation*}
q(x)=4 x_{1} x_{3}-x_{2}^{2}, \quad P_{0}=D_{x_{1}} D_{x_{3}}-D_{x_{2}}^{2} \tag{5.3}
\end{equation*}
$$

and

$$
\begin{equation*}
P=P_{0}+Q, \quad Q \in \operatorname{Diff}_{3,(1-2-3)}(\tilde{X}) \tag{5.4}
\end{equation*}
$$

The surface $\partial K$ is smooth and spacelike for $P$ away from $\Gamma$ (if $\widetilde{X}$ is sufficiently small). More precisely:


Figure 5.1. $K_{1} \cap \partial \widetilde{X}_{1}, \beta^{*} \widetilde{F} \cap \partial \widetilde{X}_{1}$ and $\beta^{*} Q_{ \pm} \cap \partial \widetilde{X}_{1}$

Proposition 5.2. - If $r(x, y)=\left(\sum_{1 \leq i \leq 3} x_{i}^{12 / i}\right)^{1 / 12}$ is the 1-2-3 radial variable and $\mathbf{S}_{z}^{-}, z \in \widetilde{X}$, is the retarded solid $P$-cone over $z$ then

1. $\mathbf{S}_{z}^{-} \cap \tilde{X} \subset K \cap \tilde{X}$ if $z \in K \cap \tilde{X}$.
2. $\inf \left\{r\left(z^{\prime}\right): z^{\prime} \in \mathbf{S}_{z}^{-}\right\} \geq C^{-1} r(z)$ if $z \in K \cap \widetilde{X}$.

Proof. - Since the singular set of $\partial K$ is of codimension 2, and thus the smooth part, $\partial K \backslash \Gamma$, is connected, it suffices for part a) to show that $\partial K \backslash \Gamma$ is space like, i.e. that for $m \in \partial K \backslash \Gamma$ and $f$, a defining function of $\partial K$ near $m, p(d f(m))>0$. As the function $f$ we can take $\beta_{*} q_{0}+\epsilon$ which is smooth in a sufficiently small neighbourhood of $m$. Thus we want

$$
p\left(d\left(\beta_{*} q_{0}\right)\right)=r^{-4} p_{\beta}\left({ }^{b} d q_{0}\right)>c r^{-4}, \quad p_{\beta}={ }^{b} \sigma_{2}\left(r^{4} \beta^{*} P \beta_{*}\right), \quad c>0
$$

From (5.4) we see that $\beta^{*} P \beta_{*}=\beta^{*} P_{0} \beta_{*}+r Q_{1}, Q_{1} \in \operatorname{Diff}_{b}^{2}\left(\widetilde{X}_{1}\right)$. Hence, by shrinking the domain, it suffices to show that

$$
p_{0 \beta}\left({ }^{b} d q_{0}\right)>c, \quad p_{0 \beta}={ }^{b} \sigma_{2}\left(r^{4} \beta^{*} P_{0} \beta_{*}\right)
$$

for $q_{0}=-\epsilon$. Since $q_{0}$ is independent of $r$ and $y$, the left hand side can computed in local coordinates on $\partial \widetilde{X}_{1}$, yielding $c$ proportional to $\epsilon$.
To establish b) let us consider the retarded cone $S_{x}^{-}, x \in \mathbb{R}^{3}$, over $\{(x, y): y \in$ $\left.\mathbb{R}^{n-3}\right\} \cap N$ where $N$ is a small neighbourhood of $\Gamma \in \mathbb{R}^{n}$. It is defined as the projection to $N$ of the union of the maximally extended retarded bicharacteristics starting at $N^{*}\left(\left\{(x, y): y \in \mathbb{R}^{n-3}\right\} \cap N\right) \cap p^{-1}(0)$. By the analogy with $\mathbf{S}_{z}^{-}$we define
the solid retarded cone over $\left\{(x, y): y \in \mathbb{R}^{n-3}\right\} \cap N, \mathbf{S}_{x}^{-}$. We then claim that for $N$ small enough

$$
\begin{equation*}
\mathbf{S}_{(x, y)}^{-} \cap N \subset \mathbf{S}_{x}^{-} \cap N \tag{5.5}
\end{equation*}
$$

for which it is enough to show that the tangent cone of $S_{(x, y)}^{-}, T_{(x, y)} S_{(x, y)}^{-}$is contained in a component of $\mathbb{R}^{n} \backslash T_{(x, y)} S_{x}^{-}$, where $T_{(x, y)} S_{x}^{-}$is the tangent cone of $S_{x}^{-}$. Since $S_{x}^{-}$is characteristic, $T_{(x, y)} S_{x}^{--}$is tangent to $T_{(x, y)} S_{(x, y)}^{-}$and the conclusion follows from the general fact about quadratic forms: let $g$ be a Lorentz quadratic form (i.e. of indices of inertia $1, n-1$ ). If $C_{ \pm}$are the connected components of the cone $C=\{X: g(X)>0\}$ and $\Pi$ is a plane tangent to $C$ then $\Pi$ separates $C_{+}$and $C_{-}$.
 neighbourhood of a point in $\Gamma$ would be contained in $\mathbf{S}_{z}^{-\circ}(\Gamma \subset \partial K)$.
Thus the minimal value of $r\left(z^{\prime}\right), z^{\prime} \in \mathbf{S}_{z}^{-}, z \in K \cap \widetilde{X}$ is attained on $\partial \mathbf{S}_{z}^{-}=S_{z}^{-}$. We also note that (5.5) implies

$$
\begin{equation*}
\inf _{z \in S_{(x, y)}^{-} \cap N} r(z) \geq \inf _{z \in S_{x}^{-} \cap N} r(z) \tag{5.6}
\end{equation*}
$$

so we would like the bound for the right hand side.
We will first obtain it for $Q_{x}^{-}$, that is, for $P=P_{0}$, the operator with constant coefficients and homogeneous of degree -4 . In that case $Q_{x}=\left\{\left(x^{\prime}, y^{\prime}\right): q\left(x^{\prime}-x\right)=0\right\}$ for a 1-2-3 homogeneous polynomial of degree $4, q$. If $r(x, y)=1$ we easily see that $0 \notin Q_{x}^{-}$. In fact, we would then have $(x, y) \in Q_{0}^{+}$which contradicts $\beta^{*} K \cap \beta^{*} Q_{0}^{+}=\emptyset$ if $(x, y) \in K$. Thus, for $(x, y) \in K, r(x, y)=1$, we have $r\left(x^{\prime}, y^{\prime}\right) \geq C^{-1}$ if $\left(x^{\prime}, y^{\prime}\right) \in Q_{x}^{-}$. For any $(x, y)=z \in K$ we obtain $(\tilde{x}, \tilde{y})=\tilde{z}=T_{r(z)^{-1}} z \in K$ with $r(\tilde{z})=1$. Since $q$ is homogeneous we then have $Q_{x}^{-}=\left(T_{r(z)}\right) Q_{\tilde{x}}^{-}$and for $z^{\prime} \in Q_{x}^{-}$,

$$
r\left(z^{\prime}\right) \geq r(z) \inf \left\{r(v): v \in Q_{\tilde{x}}^{-}\right\} \geq C^{-1} r(z)
$$

The general case follows from a perturbation argument based on (5.5) and (5.6).
Let us first observe that

$$
\inf _{z \in S_{(x, y)}^{-}} r(z)=\inf _{\gamma_{(x, y)}^{-}} \inf _{z \in \gamma_{(x, y)}^{-}} r(z)
$$

where $\gamma_{(x, y)}^{-}$runs through all the retarded characteristics starting at $(x, y)$. For $\gamma=\gamma_{(x, y)}^{-}$, let us define

$$
I_{\gamma}=\inf \{r(\gamma(t)): 0<t<\inf \{t: r(\gamma(t))>2 r(x, y)\}\}
$$

and it is enough to show that for some $c>0$ and $(x, y) \in K \cap N, I_{\gamma_{(x, y)}^{-}}>c r(x, y)$. In fact,

$$
\begin{aligned}
& \inf \left\{r(\gamma(t)): 0<t<t_{\infty}, \gamma(t) \in N\right\}=\inf _{k} \inf \left\{r(\gamma(t)): t_{k}^{1}<t<t_{k}^{2}\right\}= \\
& \inf _{k} \inf _{\tilde{\gamma}_{k}=\gamma_{\gamma\left(t_{k}^{1}\right)}^{-}} \inf \left\{r\left(\tilde{\gamma}_{k}(t)\right): 0<t<t_{k}^{2}-t_{k}^{1}\right\} \geq \inf _{k} I_{\tilde{\gamma}_{k}}>c r(x, y)
\end{aligned}
$$

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if we choose $t_{0}^{1}=0$ and

$$
t_{k}^{2}=\inf \left\{t>t_{k}^{1}: r(\gamma(t))>2 r\left(\gamma\left(t_{k}^{1}\right)\right)\right\}, \quad t_{k+1}^{1}=\sup \left\{t>t_{k}^{2}: r(\gamma(t))>r\left(\gamma\left(t_{k}^{2}\right)\right)\right\}
$$

Part a) of the proposition was used here in asserting that $\gamma\left(t_{k}^{1}\right) \in K$.
In view of (5.5) we obtain

$$
I_{\gamma_{(x, y)}^{-}} \geq \inf _{y} \inf _{\tilde{\gamma}=\tilde{\gamma}_{(x, y)}^{-}} I_{\tilde{\gamma}}
$$

where $\tilde{\gamma}$ runs through the retarded characteristics lying in $S_{x}^{-}$. Consequently, we only need to show that $I_{\tilde{\gamma}}>c r(x)$ for any such $\tilde{\gamma}$.
Let us make a change of variables

$$
z=T_{\delta^{-1}} z^{\prime}, \quad z=(x, y), z^{\prime}=\left(x^{\prime}, y^{\prime}\right)
$$

so that

$$
P\left(z^{\prime}, D_{z^{\prime}}\right)=\delta^{-4}\left(P_{0}\left(x, D_{x}\right)+\delta P_{1}\left(\delta, z, D_{z}\right)\right)
$$

and we introduce the operator $P_{\delta}\left(z, D_{z}\right)=P_{0}+\delta P_{1}$. Using its principal symbol $p_{\delta}$, we can define $\tilde{\gamma}_{(x, y), \delta}^{-}$and $S_{x, \delta}^{ \pm}$as we did for $p$. It follows that

$$
\begin{equation*}
\tilde{\gamma}_{T_{\delta}(x, y)}^{-}=T_{\delta} \tilde{\gamma}_{(x, y), \delta}^{-}, \quad S_{T_{\delta} x}^{-}=T_{\delta} S_{x, \delta}^{-}, \tag{5.7}
\end{equation*}
$$

To describe $\tilde{\gamma}_{(x, y), \delta}^{-}$we can use the parametrization

$$
\tilde{\gamma}_{(x, y), \delta}^{-}(t)=\pi\left(\exp \left(t H_{p_{\delta}}(x, y ; \xi, 0)\right)\right), \quad|\xi|=1, t>0
$$

with $(x, y ; \xi, 0)$ in the retarded component of

$$
p_{\delta}^{-1}(0) \cap N_{(x, y)}^{*}\left\{(x, \tilde{y}): \tilde{y} \in \mathbb{R}^{n-3}\right\} .
$$

Since $P_{0}\left(x, D_{x}\right)$ is strictly hyperbolic in $\mathbb{R}^{3}$ (or by a direct computation) there exists $T$ independent of $0<\delta<\delta_{0}$ such that for $T<t<2 T, r\left(\tilde{\gamma}_{(x, y), \delta}^{-}(t)\right)>2$. This is obtained by first arranging that $d(\Gamma, \tilde{\gamma}(t))>C, T<t<2 T$ where $d$ is the euclidean distance and then observing that $r(z)>C^{\prime} \min \left(d(z, \Gamma), d(z, \Gamma)^{\frac{1}{3}}\right)$. Hence for $r(x)=1$,

$$
I_{\tilde{\gamma}_{T_{\delta}(x, y)}^{-}}^{-} \geq \inf _{0<t<T} r\left(T_{\delta} \widetilde{\gamma}_{(x, y), \delta}^{-}(t)\right)
$$

and it is enough to show that there exists $\delta_{0}>0$ such that for $\delta<\delta_{0}$ and $r(x)=1$, $\inf _{0<t<T} r\left(\gamma_{(x, y), \delta}(t)\right)>c$. Since

$$
d\left(p_{\delta}^{-1}(0) \cap\{\eta=0\} \cap S^{*} N, p_{0}^{-1}(0) \cap\{\eta=0\} \cap S^{*} N\right) \longrightarrow 0, \delta \longrightarrow 0
$$

the continuous dependence on parameters and initial data for solution of ordinary differential equations implies that

$$
\sup _{0<t<T} d\left(\tilde{\gamma}_{(x, y), 0}^{-}(t), \tilde{\gamma}_{(x, y), \gamma}^{-}(t)\right) \longrightarrow 0, \delta \longrightarrow 0
$$

On the other hand $Q_{x}^{-}=S_{x, 0}^{-}$and hence

$$
\inf r\left(\tilde{\gamma}_{(x, y), 0}(t)\right) \geq \inf _{z \in Q_{x}^{-}} r(z) \geq C^{-1}
$$

By taking $\delta$ small enough we get $c=1 / 2 C$.
If $\phi(z)$ is a time function for $P$ (in view of Proposition 3.1, we can take $\phi(x, y)=$ $x_{1}+x_{3}$ ) once we arrange in Lemma 5.5 below that $\left.\sigma_{2}(Q)\right|_{N^{*} \Gamma}=0$ we can assume that $\{z:|\phi(z)|<\delta\} \subset \widetilde{X} \subset\{z:|\phi(z)|<2 \delta\}$ for some small $\delta$. We then define $\widetilde{X}_{-}=\{z \in \widetilde{X}: \phi(z)<-\delta / 2\}$ and denoting $K \cap \widetilde{X}$ by $K$,

$$
L_{-}^{2}(K)=\left\{u \in L^{2}(K): u \equiv 0 \text { in } K \cap \tilde{X}_{-}\right\}
$$

If $A$ is the forward fundamental solution of $P$ we form the operators

$$
B_{j}= \begin{cases}D_{x_{j}} A & 1 \leq j \leq 3  \tag{5.8}\\ D_{y_{j-3}} A & 4 \leq j \leq n\end{cases}
$$

which in view of part a) of Proposition 5.2 and the energy estimate, have the mapping property

$$
B_{j}: L_{-}^{2}(K) \longrightarrow L_{-}^{2}(K), \quad 1 \leq j \leq n
$$

We can now state
Lemma 5.3. - If $A$ is the forward fundamental solution for $P$ and $K, \widetilde{X}$ and $B_{j}$ are defined above then

$$
\begin{equation*}
B_{j}: r^{-p} L_{-}^{2}(K) \longrightarrow r^{-p} L_{-}^{2}(K) \tag{5.9}
\end{equation*}
$$

for all $p \geq 1$ and $1 \leq j \leq n$.
Proof. - By Proposition 5.2, $A f \in \mathscr{D}^{\prime}(K)$ is well defined for $f \in L_{-}^{2}(K)$ and the energy inequality gives (5.9) for $p=0$. If $b_{j}$ is the Schwartz kernel of $B_{j}$, then, since supp $b_{j}(z, \bullet) \subset \mathbf{S}_{z}^{-}$, part b) of Proposition 5.2 gives

$$
\begin{equation*}
b_{j}\left(z, z^{\prime}\right)=0 \quad \text { if } \quad r\left(z^{\prime}\right)<c r(z), z, z^{\prime} \in K \tag{5.10}
\end{equation*}
$$

One can now apply the dyadic decomposition argument as in [32]:

$$
K=\bigcup K_{j}^{\rho}, \quad K_{j}^{\rho}=K \cap\left\{z: 2^{-j-1} \rho \leq r(z) \leq 2^{-j} \rho\right\}, \quad K_{(j)}^{\rho}=\bigcup_{k \leq j} K_{k}^{\rho}
$$

Let us define

$$
u_{(j)}^{\rho}(z)= \begin{cases}u(z), & z \in K_{(j)}^{\rho} \\ 0 & z \in K \backslash K_{(j)}^{\rho} .\end{cases}
$$

It is easy to see that with $L^{2}(K)$ norms and $p>0$,

$$
\begin{equation*}
\left\|r^{p} u\right\|^{2} \leq \rho^{p} \sum_{j} 2^{-2 j p}\left\|u_{(j)}^{\rho}\right\|^{2} \leq \frac{2^{4 p}}{2^{2 p}-1}\left\|r^{p} u\right\|^{2} \tag{5.11}
\end{equation*}
$$

It now follows from (5.10) that $(B u)_{(j)}^{1}=\left(B\left(u_{(j)}\right)^{c}\right)_{(j)}^{1}$ and from (5.11) that

$$
\left\|r^{p} B u\right\|^{2} \leq c^{-p} \frac{2^{4 p}}{2^{2 p}-1}\left\|r^{p} u\right\|^{2}
$$

which proves the lemma.
We now show that if the only singular data for the free propagation in the past is the incident front, no new singularities appear in the past on the resolved level as well (see Fig. 5.1). In particular, the retarded cone is not present at all.

Since, away from $\Gamma, \widetilde{F}$ and $\widetilde{S}$ are simply tangent along $B$ we find that in coordinates in which (5.3) holds

$$
\widetilde{F}=\left\{(x, y): x_{3}=x_{2}^{2} a\left(x_{1}, x_{2}, y\right)\right\}, \quad a \in C^{\infty}
$$

The map

$$
\left(x_{1}, x_{2}, x_{3}, y\right) \mapsto\left(x_{1}, x_{2}\left(1-x_{1} a\right)^{\frac{1}{2}}, x_{3}-x_{2}^{2} a, y\right)
$$

preserves $\widetilde{S}, B$ and maps $\widetilde{F}$ into

$$
\widetilde{F}=\left\{(x, y): x_{3}=0\right\}
$$

It is also convenient to assume, without any loss of generality, that the coefficient of $D_{x_{2}}^{2}$ in $P$ is 1.

Proposition 5.4. - If the operators $B_{j}$ are defined by (5.8) and $K_{1}$ is given by Definition 5.1 in the definition of $K, K_{1}=\beta^{*} K$, then

$$
B_{j}: \beta_{*} I_{k} L_{\nu}^{2}\left(K_{1}, \partial \widetilde{X}_{1} \sqcup \beta^{*} \widetilde{F}\right) \cap L_{-}^{2}(K) \longrightarrow \beta_{*} I_{k} L_{\nu}^{2}\left(K_{1}, \partial \widetilde{X}_{1} \sqcup \beta^{*} \widetilde{F}\right)
$$

Before proceeding with the proof we shall establish
Lemma 5.5. - Let us define the following vector fields:

$$
V_{00}=\sum_{j=1}^{3} j x_{j} D_{x_{j}}, V_{0 j}=D_{y_{j}}, 1 \leq j \leq n-3, V_{11}=D_{x_{1}}, \quad V_{21}=D_{x_{2}}
$$

Then, after a possible change of coordinates satisfying (3.8),

$$
\left[P, V_{i k}\right]=-4 i \delta_{k 0} P+\sum_{j, l} W_{i k j l} V_{j l}+Z_{i k}
$$

where $W_{i k j l}, Z_{i k} \in \operatorname{Diff}^{1}\left(X_{1}\right)$ and the coefficients of $W_{i k j l}$ are $\mathcal{B}\left(r^{\max (0, j-i)}\right)$.

Proof. - We first want to arrange that $\left.\tau_{2}(Q)\right|_{N^{*} \Gamma}=0$ while preserving the form of $\widetilde{F}$ (with a change of coordinates satisfying (3.8)). That however is quite easy as the only terms in $Q$ with homogeneous symbols possibly non-zero on $N^{*} \Gamma$ are $a D_{x_{1}}^{2}+b D_{x_{1}} D_{x_{2}}$. Thus replace $a$ and $b$ by $a(0, y)$ and $b(0, y)$ and make the change of variables

$$
x_{1} \longmapsto x_{1}-\frac{b}{2} x_{2}+\left(a-\frac{1}{4} b^{2}\right) x_{3}, \quad x_{i} \longmapsto x_{i}, \quad i \neq 1 .
$$

Since then $D_{x_{1}} \mapsto D_{x_{1}}, D_{x_{2}}+\frac{b}{2} D_{x_{1}} \mapsto D_{x_{2}}$ and $D_{x_{3}}-\left(a-\frac{b^{2}}{4}\right) D_{x_{1}} \mapsto D_{x_{3}}$, the reduction is complete as $D_{x_{2}}^{2}-D_{x_{1}} D_{x_{3}}+a D_{x_{1}}^{2}+b D_{x_{1}} D_{x_{2}} \mapsto D_{x_{2}}^{2}-D_{x_{1}} D_{x_{3}}$.
Since $\left[P_{0}, V_{i k}\right]$ can be explicitly calculated and $\left\{x_{3}=0\right\}$ is characteristic for $P$, it is easy to see that we only need to check this for $\left[Q, V_{i k}\right], Q$ as in (5.4). Thus $W_{0 k 11} V_{11}$ contains only terms $D_{x_{1}} D_{x_{2}}, D_{x_{1}}^{2}, D_{x_{3}} D_{x_{1}}$, as the terms $D_{y_{j}} D_{x_{1}}$ can be put in $W_{0 k 0 j} V_{0 j}$. Since the symbols of $Q$ and $V_{0 k}$ vanish on $N^{*} \Gamma$, the coefficients have at least one factor of $x_{i}$ and thus are $\mathcal{O}(r)$.
The term $W_{0 k 21} V_{21}$ can only have $D_{x_{3}} D_{x_{2}}$, as $D_{x_{2}}^{2}$ does not occur ( $D_{x_{2}}^{2}$ has coefficient 1 in $P$ and $D_{x_{1}} D_{x_{2}}, D_{y_{j}} D_{x_{2}}$ can be included in $W_{0 k 11} V_{11}$ and $W_{0 k 0 j} V_{0 j}$ respectively). Since $D_{x_{3}} D_{x_{2}}$ is homogeneous of degree -5 it has to have a coefficient in $\mathcal{O}\left(r^{2}\right)$.
Finally we consider $W_{1121} V_{21}$ which has terms of homogeneity -4 or higher. Again, it has to have coefficients in $\mathcal{O}(r)$.

Proof. - Proof of Proposition 5.4 We start by characterizing the push-forward of the conormal spaces in the proposition:

$$
\begin{equation*}
g \in \beta_{*} I_{k} L_{\nu}^{2}\left(\beta^{*} K, \partial \widetilde{X}_{1} \sqcup \beta^{*} \tilde{F}\right) \cap L_{-}^{2}(K) \Longleftrightarrow \prod_{i, k} V_{i k}^{\alpha_{i k}} g \in r^{-2 \alpha_{21}-\alpha_{11}} L_{-}^{2}(K) \tag{5.12}
\end{equation*}
$$

$|\alpha| \leq k$, which follows easily from the definition of the conormal space in $\widetilde{X}_{1}$. If $u=A f, f \in \beta_{*} I_{k} L_{\nu}^{2}\left(\beta^{*} K, \partial \widetilde{X}_{1} \sqcup \beta^{*} \widetilde{F}\right) \cap L_{-}^{2}(K)$, we want $g=D_{x_{i}} u, D_{y_{j}} u$ to satisfy the condition on the right hand side of (5.12). To obtain that we shall use a more refined version of the usual 'system' argument (see [32], chapter 6). For that we will define

$$
\begin{equation*}
(u)_{p, m, l} \in L^{2}\left(\tilde{X} ; \mathbb{C}^{N}\right), \quad N=N(p, m, l) \in \mathbb{N} \tag{5.13}
\end{equation*}
$$

by an inductive procedure using

$$
\begin{aligned}
(u)_{p, m, l}^{\sharp}= & \left\{V_{0 k}(u)_{p-1, l, m}^{\sharp} ; V_{i k}(u)_{p, m, l-1}^{\sharp} ; V_{j k}(u)_{p, m-1, l}^{\sharp}:\right. \\
& i=0,1 ; j=0,1,2 ; k=0, \cdots, n-3\},(u)_{0,0,0}^{\sharp}=u,
\end{aligned}
$$

where the vector fields are applied to each component of $(u)_{0}^{\sharp}$, assumed to be 0 if any of the indices is less than 0 . Note that we allow repetitions and order the components of $(u)^{\sharp}$ using the lexicographic ordering of products of vector fields. To define (5.13)

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we start by setting $(u)_{0,0,0}=u$ and continue with

$$
\begin{aligned}
& (u)_{0,0, l+1}=\left\{(u)_{0,0, l} ;(u)_{0,0, l+1}^{\sharp}\right\}, \\
& (u)_{0, m+1, l}=\left\{(u)_{0, m, l} ;(u)_{0, m+1, l}^{\sharp}\right\}, \\
& (u)_{p+1, m, l}=\left\{(u)_{p, m, l} ;(u)_{p+1, m, l}^{\sharp} \cdot\right\}
\end{aligned}
$$

where to identify with $\mathbb{C}^{N}, N=N(p, m, l)$ we again use the lexicographic ordering of products of vector fields. We again adopt a convention that $(u)_{p, m, l}=0$ if any of the indices is less than 0.
The characterization (5.12) can be rephrased in terms of $(g)_{p, m, l}$ as follows:

$$
\begin{equation*}
g \in \beta_{*} I_{k} L_{\nu}^{2}\left(\beta^{*} K, \partial \widetilde{X}_{1} \sqcup \beta^{*} \widetilde{F}\right) \cap L_{-}^{2}(K) \Longleftrightarrow(g)_{p, m, l} \in r^{-m-2 l} L_{-}^{2}\left(K ; \mathbb{C}^{N(p, m, l)}\right), \tag{5.14}
\end{equation*}
$$

$p+m+l \leq k$. If $P u=f$, we will obtain the following system

$$
\begin{align*}
\mathscr{P}_{p, m, l}(u)_{p, m, l}= & \overbrace{p, m, l}^{20}(u)_{p-1, m, l+1}+\overbrace{p, m, l}^{10}(u)_{p-1, m+1, l}+  \tag{5.15}\\
& \overbrace{p, m, l}^{11}(u)_{p, m-1, l+1}+\mathcal{F}_{p, m, l}(f)_{m, l, p}
\end{align*}
$$

where

$$
\begin{equation*}
\mathscr{P}_{p, m, l} \in \operatorname{Diff}^{2}\left(\tilde{X} ; \mathbb{C}^{N(p, m, l)}, \mathbb{C}^{N(p, m, l)}\right), \quad \sigma_{2}\left(\mathscr{P}_{p, m, l}\right)=\sigma_{2}(P) I d_{\mathbb{C}^{N(p, m, l)}} \tag{5.16}
\end{equation*}
$$

$$
\begin{aligned}
& \overbrace{p, m, l}^{20} \in \operatorname{Diff}^{1}\left(\tilde{X} ; \mathbb{C}^{N(p, m, l)}, \mathbb{C}^{N(p-1, m, l+1)}\right), \\
& \overbrace{p, m, l}^{10} \in \operatorname{Diff}^{1}\left(\tilde{X} ; \mathbb{C}^{N(p, m, l)}, \mathbb{C}^{N(p-1, m+1, l)}\right), \\
& \overbrace{p, m, l}^{11} \in \operatorname{Diff}^{1}\left(\tilde{X} ; \mathbb{C}^{N(p, m, l)}, \mathbb{C}^{N(p, m-1, l+1)}\right),
\end{aligned}
$$

with the coefficients of $\mathscr{\vartheta}_{p, m, l}^{i j}$ in $\mathcal{O}\left(r^{i}\right)$, and $\mathcal{F}_{p, m, l} \in C^{\infty}\left(\widetilde{X} ; \mathbb{C}^{N(p, m, l)} \otimes \mathbb{C}^{N(p, m, l)}\right)$. If (5.16) holds then Lemma 5.5 is still available:

$$
\begin{align*}
& {\left[\mathscr{P}_{p, m, l}, V_{i k} I d_{\mathbb{C}^{N(p, m, l)}}\right]=-4 i \delta_{k 0} \mathscr{P}_{p, m, l}+}  \tag{5.17}\\
& \sum_{r, s} W_{i k r s} V_{r s} I d_{\mathbb{C}^{N(p, m, l)}}+\mathcal{Z}_{i k}^{p, m, l}, \mathcal{Z}_{i k}^{p, m, l} \in \operatorname{Diff}^{1}\left(\widetilde{X}, \mathbb{C}^{N(p, m, l)}, \mathbb{C}^{N(p, m, l)}\right)
\end{align*}
$$

The operator in the system is defined by successive inductions based on (5.17) and the inductive definition of $(u)_{p, m, l}$ :

$$
\begin{aligned}
\mathscr{P}_{0,0, l+1}(u)_{0,0, l+1}= & \left\{\mathscr{P}_{0,0, l}(u)_{0,0, l} ; \mathscr{P}_{0,0, l}\left(V_{i k}(u)_{0,0, l}^{\sharp}\right)-\sum_{r, s} W_{i k r s}\left(V_{r s}(u)_{0,0, l}^{\sharp}\right)-\right. \\
& \left.\mathcal{Z}_{i k}^{0,0, l}(u)_{0,0, l}^{\sharp}, i \neq 2\right\},
\end{aligned}
$$

where we identify $V_{j k}(u)_{0,0, l}^{\sharp}$ with components of $(u)_{0,0, l+1}^{\sharp}$ but use $\mathscr{P}_{0,0, l}$ on $\left\{0 ; V_{i k}(u)_{0,0, l}^{\sharp}\right\} \in \mathscr{D}^{\prime}\left(\widetilde{X} ; \mathbb{C}^{N(0,0, l)}\right)$. The terms on the right hand side of the system
equation are ${ }^{\tau} \theta_{0,0, l}^{i j} \equiv 0$ and

$$
\begin{aligned}
& \mathcal{F}_{0,0, l+1}(f)_{0,0, l+1}=\left\{\mathcal{F}_{0,0, l}(f)_{0,0, l} ;\right. \\
&\left.V_{00} \mathscr{F}_{0,0, l}^{\sharp}(f)_{0,0, l}-4 i \mathcal{F}_{0,0, l}^{\sharp}(f)_{0,0, l}, V_{i k}\left(\mathcal{F}_{0,0, l}^{\sharp}(f)_{0,0, l}\right), i>0\right\},
\end{aligned}
$$

where $\mathcal{F}_{0,0, l}^{\sharp}(f)_{0,0, l} \stackrel{\text { def }}{=} \mathscr{P}_{0,0, l}\left(\left\{0,\left(u_{0,0, l}\right)^{\sharp}\right\}\right)$.
The next induction will introduce ${ }^{*} Q_{0, p, l}^{11}$ :

$$
\begin{gathered}
\mathscr{P}_{0, m+1, l}(u)_{0, m+1, l}=\left\{\mathscr{P}_{0, m, l}(u)_{0, m, l} ; \mathscr{P}_{0, m, l}\left(V_{j k}(u)_{0, m, l}^{\sharp}\right)-\sum_{s, r \neq 2} W_{j k r s}\left(V_{r s}(u)_{0, m, l}^{\sharp}\right)-\right. \\
\mathcal{Z}_{j k}^{0, m, l}(u)_{0, m, l}^{\sharp} ; \mathscr{P}_{0, m+1, l-1}\left(V_{i k}(u)_{0, m+1, l-1}^{\sharp}\right)-\sum_{r, s} W_{i k r s}\left(V_{r s}(u)_{0, m+1, l-1}^{\sharp}\right)- \\
\left.\mathcal{Z}_{i k}^{0, m-1, l+1}(u)_{0, m+1, l-1}^{\sharp}: j=0,1 ; i=0,1,2 ; k=0, \cdots, n-3\right\},
\end{gathered}
$$

with the same convention as before. Thus $\mathscr{P}_{0, m+1, l}$ is constructed from $\mathscr{P}_{0, m, l}$ and $\mathscr{P}_{0, m+1, l-1}$. Now ${ }^{\mathscr{C}}{ }_{0, m+1, l}^{i 0} \equiv 0$ but

$$
\mathscr{\vartheta}_{0, m+1, l}^{11}(u)_{0, m, l+1}=\{\vartheta \overbrace{0, m, l}^{11}(u)_{0, m-1, l+1} ; W_{j k 21}\left(V_{21}(u)_{0, m, l}^{\sharp}\right) ; 0 ; j \neq 2\},
$$

where we consider $V_{21}(u)_{0, m, l}^{\sharp}$ as a component of $(u)_{0, m, l+1}^{\sharp}$. We also have

$$
\begin{gathered}
\mathcal{F}_{0, m+1, l}(f)_{0, m+1, l}=\left\{\mathscr{F}_{0, m, l}(f)_{0, m, l} ; V_{00}(f)_{0, m, l}^{\sharp}-4 i \mathcal{F}_{0, m, l}^{\sharp}(f)_{0, m, l}, V_{1 k} \mathscr{F}_{0, m, l}^{\sharp}(f)_{0, m, l} ;\right. \\
V_{00}\left(\mathscr{F}_{0, m+1, l-1}^{\sharp}(f)_{0, m+1, l-1}\right)-4 i\left(\mathcal{F}_{0, m+1, l-1}^{\sharp}(f)_{0, m+1, l-1}\right), \\
\left.V_{i k}\left(\mathscr{F}_{0, m+1, l-1}^{\sharp}(f)_{0, m+1, l-1}\right) ; i=1,2 ; k=0, \cdots, n-3\right\}
\end{gathered}
$$

where $\mathscr{F}_{0, m, l}^{\sharp}(f)_{0, m, l} \stackrel{\text { def }}{=} \mathscr{P}_{0, m, l}\left(\left\{0,(u)_{0, m, l}^{\sharp}\right\}\right)$.
Finally, we define

$$
\begin{gathered}
\mathscr{P}_{p+1, m, l}(u)_{p+1, m, l}=\left\{\mathscr{P}_{p, m, l}(u)_{p, m, l} \mathscr{P}_{p, m, l}\left(V_{0 k}(u)_{p, m, l}^{\sharp}\right)-\sum_{s} W_{0 k 0 s}\left(V_{0 s}(u)_{p, m, l}^{\sharp}\right)-\right. \\
\mathcal{Z}_{0 k}^{p, m, l}(u)_{p, m, l}^{\sharp} \mathscr{P}_{p+1, m-1, l}\left(V_{j k}(u)_{p+1, m-1, l}^{\sharp}\right)-\sum_{s, r \neq 2} W_{j k r s}\left(V_{r s}(u)_{p+1, m-1, l}^{\sharp}\right)- \\
\mathcal{Z}_{j k}^{p+1, m-1, l}(u)_{p+1, m-1, l}^{\sharp} ; \mathcal{P}_{p+1, m, l-1}^{\sharp}\left(V_{i k}(u)_{p+1, m, l-1}^{\sharp}\right)-\sum_{r, s} W_{i k r s}\left(V_{r s}(u)_{p+1, m, l-1}^{\sharp}\right)- \\
\left.\mathcal{Z}_{i k}^{p+1, m, l+1}(u)_{p+1, m, l-1}^{\sharp}: j=0,1 ; i=0,1,2 ; k=0, \cdots, n-3\right\},
\end{gathered}
$$

and
$\overbrace{p+1, m, l}^{10}(u)_{p, m+1, l}=\left\{\mathcal{V O}_{p, m, l}^{10}(u)_{p-1, m+1, l} ; W_{0 k 11}\left(V_{11}(u)_{p, m, l}^{\sharp}\right) ; 0 ; 0 ; k=0, \cdots, n-3\right\}$, $\overbrace{p+1, m, l}^{20}(u)_{p, m, l+1}=\{\overbrace{p, m, l}^{20}(u)_{p-1, m, l+1} ; W_{0 k 21}\left(V_{21}(u)_{p, m, l}^{\sharp}\right) ; 0 ; 0 ; k=0, \cdots, n-3\}$,

$$
\begin{aligned}
& \overbrace{p+1, m, l}^{11}(u)_{p+1, m-1, l+1}=\left\{\text { ชิ }{ }_{p, m, l}^{11}(u)_{p, m-1, l+1} ; 0 ; W_{j k 21}\left(V_{21}(u)_{p+1, m-1, l}^{\sharp}\right) ; 0 ;\right. \\
& j=0,1 ; k=0, \cdots, n-3\},
\end{aligned}
$$

where again we use the convention that, for instance, $V_{11}(u)_{p, m, l}^{\sharp}$ is a component of $(u)_{p, m+1, l}^{\sharp}$. The definition of $\mathcal{F}_{p+1, m, l}$ is analogous to those in the previous cases. Since the principal symbol of $\mathscr{P}_{p, m, l}$ is $P I d_{\mathbb{C}^{N(p, m, l)}}$, Lemma 5.3 can be applied in this setting and thus

$$
\begin{gather*}
u=A f,(f)_{0,0, l} \in r^{-2 l} L^{2}\left(K, \mathbb{C}^{N(0,0, l)}\right) \Longrightarrow \\
(g)_{0,0, l} \in r^{-2 l} L^{2}\left(K, \mathbb{C}^{N(0,0, l)}\right), g=D_{x_{i}} u, D_{y_{j}} u \tag{5.18}
\end{gather*}
$$

Using the right hand side of (5.18) and (5.15) we can now prove by induction that

$$
\begin{equation*}
(g)_{0, m, l} \in r^{-m-2 l} L^{2}\left(K, \mathbb{C}^{N(0, m, l)}\right), \quad g=D_{x_{i}} u, D_{y_{j}} u \tag{5.19}
\end{equation*}
$$

In fact, $(f)_{0, m, l} \in r^{-m-2 l} L^{2}\left(K, \mathbb{C}^{N(0, m, l)}\right)$ in view of the characterization (5.14), and

$$
\begin{aligned}
& \mathcal{O}_{0, m, l}^{01}(u)_{0, m-2, l+1} \in r^{-(m-2)-2(l+1)} L^{2}=r^{-m-2 l} L^{2}, \\
& \overbrace{0, m, l}^{11}(u)_{0, m-1, l+1} \in \mathcal{O}(r) r^{-(m-1)-2(l+1)} L^{2} \subset r^{-m-2 l} L^{2},
\end{aligned}
$$

by the induction hypothesis. Thus another application of Lemma 5.3 gives (5.19). Using that as the starting point of an induction on $p$ concludes the proof.

The immediate consequence is the following refinement of Theorem 3.9:
Theorem 5.6. - If $J_{s} L^{2}(\tilde{X}, H)$ is given by Definition 3.5 and

$$
P u=f \quad \text { in } \tilde{X},\left.\quad u\right|_{\tilde{X}_{-}}=0, \quad f \in J_{s} L_{c}^{2}(\tilde{X}, H),\left.\quad f\right|_{\tilde{X}_{-}}=0
$$

then

$$
u \in J_{s}^{1} L_{\mathrm{loc}}^{2}(\tilde{X}, H)
$$

Proof. - Let us first take $s=k \in \mathbb{N}_{0}$. Then clearly

$$
\|v\|_{J_{k} L^{2}(\tilde{X}, H)} \leq\left\|\beta_{5}^{*} v\right\|_{I_{k} L_{\nu_{5}}^{2}\left(\widetilde{X}_{5}, \mathbb{®}_{5}\right)}+2\left\|\beta^{*} v\right\|_{I_{k} L_{\nu_{1}}^{2}\left(K_{1}, \partial X \sqcup \beta^{*} F\right)} .
$$

If $v=\chi u$ or $v=D_{j} \chi u$, where $\chi \in C_{0}^{\infty}(\tilde{X})$, then by Theorem 3.9 and Proposition 5.4 the right hand side above is bounded by

$$
\left\|\beta_{5}^{*} f\right\|_{I_{k} L_{\nu_{5}}^{2}\left(\widetilde{X}_{5}, \otimes_{5}\right)}+2\left\|\beta^{*} f\right\|_{I_{k} L^{2}\left(K_{1}, \partial X \sqcup \beta^{*} F\right)} \leq 2\|f\|_{J_{k} L^{2}(\tilde{X}, H)}
$$

Thus,

$$
\|\chi u\|_{J_{k}^{1} L^{2}(\tilde{X}, H)} \leq 2\|f\|_{J_{k} L^{2}(\tilde{X}, H)}
$$

for any $\chi \in C_{0}^{\infty}(\tilde{X})$, so that the general case is immediate by interpolation.
We would now like to have an analogue of Proposition 5.4 for the Dirichlet problem. It is convenient to find appropriate coordinate functions.

Lemma 5.7. - There exists a diffeomorphism satisfying 3.8 such that the coordinates given by Proposition 3.1 to transform to coordinates ( $x, y$ ) near $\Gamma=\{x=0\}$ in which

$$
\begin{gather*}
P=P_{0}+Q, \quad P_{0}=D_{x_{2}}^{2}-x_{2} D_{x_{1}}^{2}-D_{x_{1}} D_{x_{3}}, \\
Q \in \operatorname{Diff}_{3,1-2-3}^{2}, \quad \sigma_{2}(Q) \upharpoonright_{N * \Gamma}=0 . \tag{5.20}
\end{gather*}
$$

and

$$
\begin{equation*}
X=\left\{x_{2} \geq 0\right\}, \quad \partial X=\left\{x_{2}=0\right\} \tag{5.21}
\end{equation*}
$$

Proof. - Starting with the coordinates given by Proposition 3.1 we take an inverse near 0 of the transformation

$$
\begin{align*}
& x_{1} \longmapsto x_{1}  \tag{5.22}\\
& x_{2} \longmapsto(16)^{\frac{1}{3}}\left(3 x_{2}-x_{1}^{2}\right)  \tag{5.23}\\
& x_{3} \longmapsto 27 x_{3}+9 x_{1} x_{2}+5 x_{1}^{3} \tag{5.24}
\end{align*}
$$

so that the cusp is given by $\left\{\left(27 x_{3}+9 x_{1} x_{2}+5 x_{1}^{3}\right)^{2}-16\left(3 x_{2}-x_{1}^{2}\right)^{3}=0\right\}$. The operator for which the conormal bundle of this cusp is characteristic has to be of the form $P_{0}+Q, Q \in \operatorname{Diff}_{3,1-2-3}^{2}\left(\mathbb{R}^{n}\right)$ where we rescale $x_{1}$ if necessary - this can be seen directly or by transforming the statement of the proof in Proposition 3.1 (the motivation for this computation and the ones below is given by the explicit expressions for $\widetilde{\Lambda}_{m}$ in chapter 7).
As in the proof of Proposition 3.1 we can preserve the cusp and take the boundary to be $\left\{x_{2}=0\right\}$. Then also

$$
\widetilde{F}=\left\{(x, y): x_{3}-x_{1} x_{2}+\frac{1}{3} x_{1}^{3}+f_{4}=0\right\}, \quad f_{4} \in 9 \prod_{4}^{1-2-3}\left(\mathbb{R}^{n}\right)
$$

We claim that we can preserve $\partial X=\left\{X_{2}=0\right\}$ and map $\widetilde{F}$ to $\left\{x_{3}-x_{1} x_{2}+\frac{1}{3} x_{1}^{3}=0\right\}$ by a transformation satisfying (3.8). In fact, proceeding as in the proof of iii) in Proposition 3.3 we need a family of vector fields $V_{s}$ tangent to $\partial X$ such that

$$
V_{s}\left(x_{3}-x_{2} x_{1}+\frac{1}{3} x_{1}^{3}+s f_{4}\right)=-f_{4}
$$

and for that we can simply take, near $(0, y)$,

$$
V_{s}=\frac{-f_{4}}{1+s \partial_{x_{3}} f_{4}} \partial_{x_{3}}, \quad \partial_{x_{3}} f_{4}(0, y)=0
$$

The resulting diffeomorphism clearly satisfies (3.8). It remains to check that

$$
\sigma_{2}(Q) \upharpoonright_{N^{*} \Gamma}=0
$$

Since $\left\{x_{3}-x_{1} x_{2}+\frac{1}{3} x_{1}^{3}=0\right\}$ is characteristic for $P$ and $P_{0}$ it has to be characteristic for $Q$. The only terms with symbols nonvanishing on $N^{*} \Gamma$ are $a D_{x_{1}}^{2}+b D_{x_{1}} D_{x_{2}}$ so that leads to the condition $a\left(x_{1}^{2}-x_{2}\right)=b x_{1}$ implying $b(0, y)=a(0, y)=0$. This concludes the proof of the lemma.

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We shall need a number of preliminary facts and we start by considering the following problem. Let $V$ be a fixed vector field of the form

$$
V=a D_{x_{3}}+b D_{x_{1}}+\left\langle c, D_{y}\right\rangle, \quad a, b \in \bar{C}^{\infty}(X, \mathbb{R}), c \in \bar{C}^{\infty}\left(X, \mathbb{R}^{n-3}\right)
$$

If $\phi$ denotes a time function, let us consider the solution of

$$
\begin{equation*}
P u=V f,\left.\quad u\right|_{\partial X}=g, \quad u \equiv 0 \text { for } \phi(z)<-\delta \tag{5.25}
\end{equation*}
$$

where $f \in C_{0}^{\infty}\left(\bar{X}_{\delta}\right), g \in C_{0}^{\infty}\left(\bar{X}_{\delta} \cap \partial X\right)$ with $X_{\delta}$, a neighbourhood of 0 in $X$ such that $\phi(z)>-\delta$ in $\bar{X}_{\delta}$. We then define the map

$$
T:(f, g) \longmapsto u
$$

allowing also the notation $T_{1} g=T(0, g)$ and $T_{2} f=T(f, 0)$ so that $T(f, g)=$ $T_{1} g+T_{2} f$. Before proceeding with the analogue of Lemma 5.3 for $T_{1}, T_{2}$ we need the following

Lemma 5.8. - If $f \in C_{0}^{\infty}\left(\bar{X}_{\delta}\right)$ and $g \in C_{0}^{\infty}\left(\bar{X}_{\delta} \cap \partial X\right)$ with $X_{\delta}$ sufficiently small, bicharacteristically convex neighbourhood of 0 in $X$, then the solution of (5.25) satisfies

$$
\begin{equation*}
\|u\|_{L^{2}\left(X_{\delta}\right)} \leq C \delta^{2}\|f\|_{L^{2}\left(X_{\delta}\right)}+C \delta\|g\|_{L^{2}\left(\partial X \cap \bar{X}_{\delta}\right)} \tag{5.26}
\end{equation*}
$$

Proof. - By translating $0 \in X$ we shall assume that

$$
X_{\delta} \subset X_{\delta}^{\sharp}=\{z \in X:-\delta<\phi(z)<0\},
$$

where, as we may, we take $\phi(x, y)=x_{3}$, (only the tangency of the time function to the boundary is important). Let us recall the energy inequality, [14], (24.1.4), (compare also (24.1.6) there):

$$
\begin{align*}
\delta^{-1} \int_{Q}\left(\left|h^{\prime}\right|^{2}+|h|^{2}\right) d v & \leq C \delta \int_{Q}|P h|^{2} d v+ \\
& C \int_{\bar{Q} \cap \partial X}\left(|h|^{2}+\left|D_{x_{1}} h\right|^{2}+\left|D_{x_{3}} h\right|^{2}+\sum_{i=1}^{n-3}\left|D_{y_{i}} h\right|^{2}\right) d S \tag{5.27}
\end{align*}
$$

where $Q=\left\{z: x_{2}>0, x_{3} \leq 0\right\}$ and $h \in \mathcal{S}\left(\mathbb{R}^{n}\right), h=0$ if $\phi<-\delta$. Let $\chi \in C_{0}^{\infty}\left(\mathbb{R}^{n}\right)$ be supported in a sufficiently large ball around 0 , so that $\chi(z) u(z)=u(z)$ if $-\delta<\phi(z)<0$. The existence of $\chi$ follows from the finite speed of propagation (see [8], Sect.VII.8). Consequently, it suffices to establish (5.26) with $u$ replaced by $\chi u$.
We now follow, in a slightly modified way, the proof of Lemma 24.1 .5 in [14] by first introducing

$$
E_{s}\left(D^{\prime}\right)^{*}=\left(\left(1+D_{x_{1}}^{2}+\sum_{i=1}^{n-3} D_{y_{j}}^{2}\right)^{1 / 2}+i D_{x_{3}}\right)^{s}
$$

and then applying (5.27) to $h=E_{-1}\left(D^{\prime}\right)^{*} \chi u$. It is important here that $h=0$ in $\phi(z)<-\delta$, since $u=0$ there. We also recall that by Theorem B.2.4 of [14],

$$
\begin{aligned}
&\left\|E_{1}^{*}\left(D^{\prime}\right) h\right\|_{L^{2}\left(X_{\delta}^{\sharp}\right)} \leq C\left(\left\|D_{x_{1}} h\right\|_{L^{2}\left(X_{\delta}^{\sharp}\right)}^{2}+\left\|D_{x_{3}} h\right\|_{L^{2}\left(X_{\delta}^{\sharp}\right)}^{2}\right. \\
&\left.\quad+\sum_{i=1}^{n-3}\left\|D_{y_{i}} h\right\|_{L^{2}\left(X_{\delta}^{\sharp}\right)}^{2}+\|h\|_{L^{2}\left(X_{\delta}^{\sharp}\right)}^{2}\right)^{\frac{1}{2}}
\end{aligned}
$$

and
$E_{-1}\left(D^{\prime}\right)^{*}\left[P, E_{1}^{*}\right]=R_{x_{1}}\left(z, D_{z}^{\prime}\right) D_{x_{1}}+R_{x_{3}}\left(z, D_{z}^{\prime}\right) D_{x_{3}}+\sum_{i=1}^{n-3} R_{i}\left(z, D_{z}^{\prime}\right) D_{y_{i}}+R_{0}\left(z, D_{z}^{\prime}\right)$,
where $R_{x_{j}}, R_{i} \in S^{0}\left(\mathbb{R}^{n}, \mathbb{R}^{n-1}\right), D_{z}^{\prime}=\left(D_{x_{1}}, D_{x_{3}}, D_{y}\right)$, and $R_{\bullet}\left(z, D_{z}^{\prime}\right) D_{\bullet} w=$ $R_{\bullet}\left(z, D_{z}^{\prime}\right) w_{\bullet}$, if $w_{\bullet}=D_{\bullet} w$ for $\phi(v)<0$ and $w_{\bullet}=0$ elsewhere.
Thus,

$$
\begin{aligned}
\|\chi u\|_{L^{2}\left(X_{\delta}^{\sharp}\right)} & =\left\|E_{1}\left(D^{\prime}\right)^{*} h\right\|_{L^{2}\left(X_{\delta}^{\sharp}\right)} \\
& \leq C \delta^{2}\left\|E_{-1}\left(D^{\prime}\right)^{*} P u\right\|_{L^{2}\left(X_{\delta}^{\sharp}\right)}+C \delta\left\|E_{-1}\left(D^{\prime}\right)^{*} g\right\|_{H_{(1)}\left(\bar{X}_{\delta}^{\sharp} \cap \partial X\right)},
\end{aligned}
$$

since $\chi g=g$. We conclude the proof by observing that

$$
\begin{equation*}
\left\|E_{-1}\left(D^{\prime}\right)^{*} P u\right\|_{L^{2}\left(X_{\delta}^{\sharp}\right)} \leq C\|P u\|_{H_{(0,-1)}\left(X_{\delta}\right)} \leq C\|f\|_{L^{2}\left(X_{\delta}\right)} \tag{5.28}
\end{equation*}
$$

and

$$
\left\|E_{-1}\left(D^{\prime}\right)^{*} g\right\|_{H_{(1)}\left(\bar{X}_{\delta}^{\mathrm{A}} \cap \partial X\right)} \leq\left\|E_{-1}\left(D^{\prime}\right)^{*} g\right\|_{H_{(1)}(\partial X)} \leq\|g\|_{L^{2}(\partial X)}=\|g\|_{L^{2}\left(X_{\delta} \cap \partial X\right)} .
$$

By following the proof of Lemma 24.1.6 in [14] and using (5.28) we also obtain the mapping property:

$$
\begin{equation*}
T: L^{2}\left(X_{\delta}\right) \times L^{2}\left(\bar{X}_{\delta} \cap \partial X\right) \longrightarrow L^{2}\left(X_{\delta}\right) \tag{5.29}
\end{equation*}
$$

with small norms $\mathcal{O}\left(\delta^{2}, \delta\right)$, in the first and second factors, respectively.
Let us now modify the previous notation to the boundary value problem case. Thus we consider

$$
K_{\delta}=K \cap X_{\delta}
$$

where $K$ was defined by (5.2) and $X_{\delta}$ was given in Lemma 5.8. Similarly, we define

$$
K_{\delta}^{\partial}=K_{\delta} \cap \partial X
$$

and $L_{-}^{2}\left(K_{\delta}\right), L_{-}^{2}\left(K_{\delta}^{\partial}\right)$, by analogy with $L_{-}^{2}(K)$. With this notation we have

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Lemma 5.9. - Let $K_{\delta}, K_{\delta}^{\partial}$ be as above. If $g \in L_{-}^{2}\left(K_{\delta}^{\partial}\right)$ then

$$
\begin{equation*}
\left\|r^{p} T_{1} g\right\|_{L^{2}\left(K_{\delta}\right)} \leq C \delta\left\|r^{p} g\right\|_{L^{2}\left(K_{\delta}^{\partial}\right)} \tag{5.30}
\end{equation*}
$$

and if $f \in L_{-}^{2}\left(K_{\delta}\right)$ then

$$
\begin{equation*}
\left\|r^{p} T_{2} f\right\|_{L^{2}\left(K_{\delta}\right)} \leq C \delta^{2}\left\|r^{p} f\right\|_{L^{2}\left(K_{\delta}\right)} \tag{5.31}
\end{equation*}
$$

where, $T(f, g)=T_{2} f+T_{1} g$ is the solution operator for (5.25).
Proof. - The support property required in the proof of Lemma 5.3 holds for the mixed problems as well (see [8], Sect.VII.8). Since the existence is guaranteed for $f$ and $g$ in $L^{2}$ we can again proceed using the dyadic decomposition. The details of the proof are the same.

As in the free case we shall now consider the resolved space $X_{1}=\operatorname{cl}\left(\beta^{*} X\right)$ which is a manifold with a codimension 2 corner $\partial \tilde{X}_{1} \cap \beta^{*} \partial X$. Thus we have

$$
\beta: X_{1} \longrightarrow X, \quad \partial X_{1}=\beta^{*} \partial X \cap\left(\partial \widetilde{X}_{1} \cap X_{1}\right)
$$

We also define (compare chapter 7 below)

$$
\beta^{\partial}=\left.\beta\right|_{\beta^{*} \partial X}, \quad \beta^{\partial}: \beta^{*} \partial X \longrightarrow \partial X
$$

The following estimate will be useful later:
Lemma 5.10. - Let $N \subset X$ be such that $\beta^{*} N \subset X_{1}$ is open with smooth boundary and let $u \in C^{\infty}(\widetilde{X})$ satisfy $P u=0$ in $N$. Then

$$
\begin{equation*}
\sum_{|\alpha|+|\beta| \leq l}\left\|r^{3 \alpha_{3}+2 \alpha_{2}+\alpha_{1}} D_{x}^{\alpha} D_{y}^{\beta} u\right\|_{L^{2}(N)} \leq C_{l} \sum_{|\alpha|+|\beta| \leq l}\left\|r^{3 \alpha_{3}+\alpha_{1}} D_{x_{1}}^{\alpha_{1}} D_{x_{3}}^{\alpha_{3}} D_{y}^{\beta} u\right\|_{L^{2}(N)} \tag{5.32}
\end{equation*}
$$

for any $l \geq 0$.
Proof. - Let us define

$$
\begin{aligned}
H_{(0, l)}^{b}\left(X_{1}\right)= & \left\{w \in L_{\nu}^{2}\left(X_{1}\right): w=w_{1}+w_{2}, \operatorname{supp} w_{2} \cap \beta_{1}^{*} \partial X=\emptyset, w_{2} \in H_{(l)}^{b}\left(\tilde{X}_{1}\right)\right. \\
& w_{1} \text { supported near } \beta_{1}^{*} \partial X, W^{\gamma_{1}}\left(r D_{r}\right)^{\gamma_{2}} D_{y}^{\gamma^{\prime}} \in L_{\nu}^{2}\left(X_{1}\right) \\
& \left.|\gamma| \leq l, W=r^{i} \beta_{1}^{*} D_{x_{i}} \beta_{*}, i=1,3\right\}
\end{aligned}
$$

which is simply the mixed $b$-Sobolev space based on $\tilde{X}_{1} \supset X_{1}$ and $\beta^{*} \partial X$. Clearly, $X_{1}$ can be replaced by $\beta^{*} N$ in this definition. Thus (5.32) can be rewritten as

$$
\left\|\beta^{*} u\right\|_{H_{(l)}^{b}\left(\beta^{*} N\right)} \leq C_{l}\left\|\beta^{*} u\right\|_{H_{(0, l)}^{b}\left(\beta^{*} N\right)}, \quad P_{1} \beta_{1}^{*} u=0 \text { in } \beta^{*} N
$$

where $P_{1}=r^{4} \beta^{*} P \beta_{*}$. Since $P_{1}$ is $b$-noncharacteristic with respect to $\beta^{*} \partial X$, the estimate follows from an easy modification of the proof of Theorem B.2.9 of [14]- see Appendix B.

## SEMI-LINEAR DIFFRACTION OF CONORMAL WAVES

We also need a modification of Lemma 5.5 to help us in this setting:
Lemma 5.11. - Let us define the following vector fields:

$$
V_{i 1}=D_{x_{i}}, i=1,2,3, \quad V_{0 j}=D_{y_{j}}, j=1, \cdots, n-3 .
$$

If $P$ and $X$ are of the form (5.20) and (5.21) and the coefficient of $D_{x_{2}}^{2}$ in $P$ is 1 , then for $i \neq 2$

$$
\begin{aligned}
& {\left[P, V_{0 k}\right]=\sum_{j \neq 2,3} \widetilde{W}_{0 k j l} V_{j l}+\left(a_{0 k} V_{13}+b_{0 k} D_{x_{2}}\right) V_{13}+\widetilde{Z}_{0 k}, a_{0 k}=\mathfrak{O}\left(r^{3}\right), b_{0 k}=\mathcal{O}\left(r^{2}\right),} \\
& {\left[P, V_{11}\right]=\sum_{j \neq 2,3} \widetilde{W}_{11 j l} V_{j l}+\left(a_{11} V_{13}+b_{11} D_{x_{2}}\right) V_{13}+\widetilde{Z}_{11}, a_{1 k}=\mathcal{O}\left(r^{2}\right), b_{0 k}=\mathcal{O}(r),} \\
& {\left[P, V_{31}\right]=\sum_{j \neq 2} \widetilde{W}_{31 j l} V_{j l}+\widetilde{Z}_{31},}
\end{aligned}
$$

where $\widetilde{W}_{i k j l}, \widetilde{Z}_{i k} \in \operatorname{Diff}{ }^{1}(\widetilde{X})$ and the coefficients of $\widetilde{W}_{0 k 11}$ are in $\mathcal{O}(r)$.
Proof. - The argument used in the proof of Lemma 5.5 together with Lemma 5.7 easily yields the desired statement.

The next lemma provides the crucial a priori estimate:
Lemma 5.12. - Let $K_{\delta}, K_{\delta}^{\partial}$ be as in Lemma 5.9 and let $f \in C_{0}^{\infty}(\partial X)$ satisfy $\left.f\right|_{K_{\delta}^{\partial}} \in L_{-}^{2}\left(K_{\delta}^{\partial}\right)$. If $u$ is the solution of

$$
P u=0 \quad \text { in } X,\left.u\right|_{\phi(z)<-\delta}=0,\left.u\right|_{\partial X}=f
$$

then, for $\delta$ sufficiently small

$$
\begin{equation*}
\left\|\beta^{*} u\right\|_{I_{k} L_{\nu}^{2}\left(\beta^{*} K_{\delta} ; \partial \widetilde{X}_{1} \cap \beta^{*} K_{\delta}\right)} \leq C\left\|\beta_{\partial^{*}} f\right\|_{I_{k} L_{\nu_{\partial}}^{2}\left(\beta^{*} K_{\delta}^{\partial} ; \partial \tilde{X}_{1} \cap \beta^{*} K_{\delta}^{\partial}\right)} . \tag{5.33}
\end{equation*}
$$

Proof. - To the extent that it is possible we shall follow the proof of Proposition 5.4. Thus, the boundedness of the norm on the right hand side of (5.33) is equivalent to

$$
\begin{equation*}
\prod_{i, k, i \neq 2} V_{i k}^{\alpha_{i k}} f \in r^{-3 \alpha_{31}-\alpha_{11}} L_{-}^{2}\left(K_{\delta}^{\partial}\right),|\alpha| \leq k \tag{5.34}
\end{equation*}
$$

The required estimate for the solution now takes the form:

$$
\begin{equation*}
\prod_{i, k} V_{i k}^{\alpha_{i k}} u \in r^{-3 \alpha_{31}-2 \alpha_{21}-\alpha_{11}} L_{-}^{2}\left(K_{\delta}\right),|\alpha| \leq k . \tag{5.35}
\end{equation*}
$$

To obtain the system we define $(u)_{p, m, l}$ as before using $V_{j k}, j \neq 2$ and introduce also

$$
(u)_{p, m, l, 1} \in L^{2}\left(\tilde{X} ; \mathbb{C}^{N(p, m, l)}\right),(u)_{p, m, l, 1}=D_{x_{2}}(u)_{p, m, l},
$$

with $D_{x_{2}}$ applied to all the components of $(u)_{p, m, l}$ and the order preserved, so that, for instance,

$$
(u)_{p+1, m, l, 1}=\left\{(u)_{p, m, l, 1} ; D_{x_{2}}(u)_{p+1, m, l}^{\sharp}\right\} .
$$

As in the proof of Proposition 5.4 we obtain the following system, where now Lemma 5.11 is used:

$$
\begin{aligned}
\mathscr{P}_{p, m, l}(u)_{p, m, l}= & \mathfrak{Q}_{p, m, l}^{30}(u)_{p-1, m, l+1}+\mathfrak{Q}_{p, m, l}^{10}(u)_{p-1, m+1, l}+\mathfrak{Q}_{p, m, l}^{21}(u)_{p, m-1, l+1} \\
& \mathscr{B}_{p, m, l}^{20}(u)_{p-1, m+1, l, 1}+\mathscr{B}_{p, m, l}^{11}(u)_{p, m-1, l, 1}
\end{aligned}
$$

with boundary and initial conditions

$$
\left.(u)_{p, m, l}\right|_{\phi(z)<-\delta}=0,\left.\quad(u)_{p, m, l}\right|_{\partial X}=(f)_{p, m, l}
$$

and where

$$
\mathscr{P}_{p, m, l} \in \operatorname{Diff}^{2}\left(\tilde{X} ; \mathbb{C}^{N(p, m, l)}, \mathbb{C}^{N(p, m, l)}\right), \quad \sigma_{2}\left(\mathscr{P}_{p, m, l}\right)=\sigma_{2}(P) I d_{\mathbb{C}^{N(p, m, l)}}
$$

and $\mathfrak{C}_{p, m, l}^{i j}\left(x, D_{x_{1}}, D_{x_{3}}\right), \mathscr{B}_{p, m, l}^{i j}\left(x, D_{x_{1}}, D_{x_{3}}\right) \in \operatorname{Diff}^{1}(\tilde{X})$ with coefficients in $\mathcal{O}\left(r^{i}\right)$. Since the construction is analogous to that in Proposition 5.4 we will only describe the last, most involved, inductive definition:

$$
\begin{aligned}
& \mathscr{P}_{p+1, m, l}(u)_{p+1, m, l}=\left\{\mathscr{P}_{p, m, l}(u)_{p, m, l} ; \mathscr{P}_{p, m, l}\left(V_{0 k}(u)_{p, m, l}^{\sharp}\right)\right. \\
& \quad-\sum_{s} \widetilde{W}_{0 k 0 s}\left(V_{0 s}(u)_{p, m, l}^{\sharp}\right) \\
& \quad-\widetilde{\mathcal{Z}}_{0 k}^{p, m, l}(u)_{p, m, l}^{\sharp} ; \mathscr{P}_{p+1, m-1, l}\left(V_{j k}(u)_{p+1, m-1, l}^{\sharp}\right) \\
&-\sum_{s, r \neq 2,3} \widetilde{W}_{j k r s}\left(V_{r s}(u)_{p+1, m-1, l}^{\sharp}\right) \\
& \quad-\widetilde{\mathcal{Z}}_{j k}^{p+1, m-1, l}(u)_{p+1, m-1, l}^{\sharp} ; \mathscr{P}_{p+1, m, l-1}\left(V_{i k}(u)_{p+1, m, l-1}^{\sharp}\right) \\
&-\sum_{s, r \neq 2} \widetilde{W}_{i k r s}\left(V_{r s}(u)_{p+1, m, l-1}^{\sharp}\right) \\
&\left.-\widetilde{\mathcal{Z}}_{i k}^{p+1, m, l+1}(u)_{p+1, m, l-1}^{\sharp}: j=0,1 ; i=0,1,3 ; k=0, \cdots, n-3\right\},
\end{aligned}
$$

with

$$
\begin{aligned}
& \mathscr{B}_{p+1, m, l}^{11}(u)_{p+1, m-1, l, 1}= \\
& \left\{\mathcal{B}_{p, m-1, l, 1}^{11}(u)_{p, m-1, l, 1} ; 0 ; b_{0 k} V_{31}(u)_{p+1, m-1, l, 1}^{\sharp}, b_{11} V_{31}(u)_{p+1, m-1, l, 1}^{\sharp} ; 0\right\}, \\
& \mathscr{B}_{p+1, m, l}^{20}(u)_{p, m, l, 1}=\left\{\mathscr{B}_{p, m, l}^{20}(u)_{p-1, m, l, 1} ; b_{0 k} V_{13}(u)_{p, m, l, 1}^{\sharp} ; 0 ; 0\right\}, \\
& \mathfrak{Q}_{p+1, m, l}^{30}(u)_{p, m, l+1}=\left\{\mathfrak{Q}_{p, m, l}^{30} ; a_{0 k} V_{13}\left(V_{13}(u)_{p, m, l}^{\sharp}\right) ; 0 ; 0\right\}, \\
& \mathfrak{Q}_{p+1, m, l}^{10}(u)_{p, m+1, l}=\left\{\mathfrak{Q}_{p, m, l}^{20}(u)_{p-1, m+1, l} ; \widetilde{W}_{0 k 11}\left(V_{11}(u)_{p, m, l}^{\sharp}\right) ; 0 ; 0\right\}, \\
& \mathfrak{Q}_{p+1, m, l}^{21}(u)_{p+1, m-1, l+1}= \\
& \left\{\mathfrak{Q}_{p, m, l}^{21}(u)_{p, m-1, l+1} ; 0 ; a_{0 k} V_{31}\left(V_{31}(u)_{p, m-1, l}^{\sharp}\right), a_{11} V_{31}\left(V_{31}(u)_{p, m-1, l}^{\sharp}\right) ; 0\right\},
\end{aligned}
$$

We now use the estimates (5.30) and (5.31) of Lemma 5.9 and proceed by induction starting with

$$
(u)_{0,0, l} \in r^{-3 l} L_{-}^{2}\left(K_{\delta} ; \mathbb{C}^{N(p, m, l)}\right)
$$

Since $u \in C^{\infty}(X)$, we can assume that the norms of $(u)_{p, m, l, 1}$ in $r^{-q} L_{-}^{2}\left(K_{\delta} ; \mathbb{C}^{N(p, m, l)}\right)$ are bounded and consequently we obtain

$$
\begin{align*}
\left\|r^{3 l+m}(u)_{p, m, l}\right\|_{L_{-}^{2}\left(K_{\delta}\right)} \leq & C \delta^{2}\left(\left\|r^{3 l+m}(u)_{p-1, m, l, 1}\right\|_{L_{-}^{2}\left(K_{\delta}\right)}\right.  \tag{5.36}\\
& +\left\|r^{3 l+m}(u)_{p, m_{1}, l, 1}\right\|_{L_{-}^{2}\left(K_{\delta}\right)} \\
& \left.+\left\|r^{3 l+m}(u)_{p, m-2, l, 1}\right\|_{L_{-}^{2}\left(K_{\delta}\right)}\right) \\
& +C \delta\left\|r^{3 l+m}(f)_{p, m, l}\right\|_{L_{-}^{2}\left(K_{\delta}^{\partial}\right)} .
\end{align*}
$$

Summing (5.36) in $p, m, l, p+l+m \leq k$ we obtain

$$
\begin{aligned}
& \sum_{|\alpha| \leq k}\left\|r^{3 \alpha_{31}+\alpha_{11}} D_{x_{1}}^{\alpha_{11}} D_{x_{3}}^{\alpha_{31}} D_{y}^{\alpha^{\prime}} u\right\|_{L_{-}^{2}\left(K_{\delta}\right)} \leq \\
& C \delta^{2} \sum_{|\alpha|+|\beta| \leq k}\left\|r^{3 \alpha_{31}+2 \alpha_{21}+\alpha_{11}} D_{x}^{\alpha} D_{y}^{\beta} u\right\|_{L_{-}^{2}\left(K_{\delta}\right)} \\
&+C \delta \sum_{|\alpha| \leq k}\left\|r^{3 \alpha_{31}+\alpha_{11}} D_{x_{1}}^{\alpha_{11}} D_{x_{3}}^{\alpha_{31}} D_{y}^{\alpha^{\prime}} f\right\|_{L_{-}^{2}\left(K_{\delta}^{\sigma}\right)}
\end{aligned}
$$

If $\delta$ is small enough we can apply Lemma 5.10 with $N=K_{\delta}$ to obtain the desired $a$ priori estimate.

For $f \in L_{-}^{2}(\partial X)$ we consider the Poisson operator $T_{1}$,

$$
P T_{1} f=0 \text { in } X,\left.\quad T_{1} f\right|_{\partial X}=f,\left.\quad T_{1} f\right|_{\phi(z)<-\delta} \equiv 0
$$

We can use Lemma 5.12 to deduce
Proposition 5.13. - The Poisson operator $T_{1}$ given above has the mapping property:

$$
T_{1}: \beta_{*}^{\partial} I_{k} L_{\nu_{\partial}}^{2}\left(\beta^{\partial^{*}} K_{\delta}^{\partial} ; \partial \widetilde{X}_{1} \cap \beta^{\partial^{*}} K_{\delta}^{\partial}\right) \cap L_{-}^{2}\left(K_{\delta}^{\partial}\right) \longrightarrow \beta_{*} I_{k} L_{\nu}^{2}\left(\beta^{*} K_{\delta} ; \partial \widetilde{X}_{1} \cap \beta^{*} K_{\delta}\right)
$$

where $\beta_{*} \nu=d x d y, \beta_{*}^{\partial} \nu_{\partial}=\imath^{*} d x d y, \imath: \partial X \hookrightarrow X$.
Proof. - By the comment following the proof of Lemma 5.8, the map

$$
T_{1}: L_{-}^{2}\left(K_{\delta}^{\partial}\right) \longrightarrow \mathscr{D}^{\prime}\left(K_{\delta}\right)
$$

is well defined and, since we know the existence in $L^{2}$, we have obtained $u \in L_{-}^{2}\left(K_{\delta}\right)$ such that

$$
u=\left.T_{1} f\right|_{K_{\delta}} \quad \text { if }\left.f\right|_{K_{\delta}^{\partial}} \in L_{-}^{2}\left(K_{\delta}^{\partial}\right)
$$

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The conormal spaces in the statement of the proposition are complete (as they can be identified with $H_{(k)}^{b}$ of a manifold with corners - see Appendix B) and $\left.\beta^{*} C^{\infty}(X)\right|_{\beta^{*} N_{\delta}}$, $\left.\beta_{1}^{\partial^{*}} C^{\infty}(\partial X)\right|_{\beta_{1}^{\partial *} K_{\delta}^{\partial}}$ are dense in

$$
I_{k} L_{\nu}^{2}\left(\beta^{*} K_{\delta} ; \partial \widetilde{X}_{1} \cap \beta^{*} K_{\delta}\right) \text { and } I_{k} L_{\nu_{\partial}}^{2}\left(\beta^{\partial^{*}} K_{\delta}^{\partial} ; \partial \widetilde{X}_{1} \cap \beta^{\partial^{*}} K_{\delta}^{\partial}\right) \cap L_{-}^{2}\left(K_{\delta}^{\partial}\right)
$$

respectively. This can be seen using the density of $\beta^{*} C_{0,0}^{\infty}(X)$, where $C_{0,0}^{\infty}$ are all functions vanishing to infinite order at $\Gamma=\left\{(0, y): y \in \mathbb{R}^{n-3}\right\}$ which follows from the density of $\mathfrak{S}\left(\mathbb{R}^{n}\right)$ in $H_{(k)}\left(\mathbb{R}^{n}\right)$ (where we again use the $b$-Sobolev spaces on manifolds with corners and the identification with the usual Sobolev spaces through a logarithmic change of variables).
Thus, if $F \in I_{k} L_{\nu_{\partial}}^{2}\left(\beta^{\partial^{*}} K_{\delta}^{\partial} ; \partial \widetilde{X}_{1} \cap \beta^{\partial^{*}} K_{\delta}^{\partial}\right) \cap L_{-}^{2}\left(K_{\delta}^{\partial}\right), f \in \beta_{*} F$, then there exist $\left.F_{N} \in \beta^{\partial^{*}} C^{\infty}(\partial X)\right|_{K_{\delta}^{\partial}} ^{\partial}$ such that

$$
F_{N} \longrightarrow F \text { in } I_{k} L_{\nu_{\partial}}^{2}\left(\beta^{\partial^{*}} K_{\delta}^{\partial} ; \partial \widetilde{X}_{1} \cap \beta^{\partial^{*}} K_{\delta}^{\partial}\right) \cap L_{-}^{2}\left(K_{\delta}^{\partial}\right), \quad N \longrightarrow \infty
$$

We then consider $u_{N}=T_{1} \beta^{\partial}{ }_{*} F_{N}$, where using the a priori estimate of Lemma 5.12, we conclude that $u \in I_{k} L_{\nu}^{2}\left(\beta^{*} K_{\delta} ; \partial \widetilde{X}_{1} \cap \beta^{*} K_{\delta}\right)$.

## 6. THE EXTENSION PROPERTY

The purpose of this section is to construct an extension map from the marked Lagrangian spaces in $X$ defined in Section 4 into the conormal space $\left(\beta_{5}\right)_{*}\left(I_{k} L_{\nu_{5}}^{2}\left(\widetilde{X}_{5}, \delta_{5}\right)\right)$. Let us first fix our notation (see chapter 2):

$$
\begin{equation*}
\Lambda_{R}=N^{*} \widetilde{R} \backslash 0, \quad \Lambda_{F}=N^{*} \widetilde{F} \backslash 0, \quad \Lambda_{S}=N^{*} \widetilde{S} \backslash 0, \quad \Lambda_{H}=N^{*} H \backslash 0 \tag{6.1}
\end{equation*}
$$

The main result is

Theorem 6.1. - There exist linear and continuous maps

$$
\begin{gather*}
E_{1}: I_{k}^{b} L_{c}^{2}\left(X, \Lambda_{S} \sqcup \Lambda_{H}\right)+I_{k}^{b} L_{c}^{2}\left(X, \Lambda_{S} ; \Lambda_{S} \cap \Lambda_{F}\right) \longrightarrow\left(\beta_{5}\right)_{*}\left(I_{k} L_{\nu_{5}}^{2}\left(\tilde{X}_{5}, \mathcal{S}_{5}\right)\right), \\
\text { and }  \tag{6.2}\\
E_{2}: I_{k}^{b} L_{c}^{2}\left(X, \Lambda_{F} \sqcup \Lambda_{R}\right) \longrightarrow\left(\beta_{5}\right)_{*}\left(I_{k} L_{\nu_{5}}^{2}\left(\widetilde{X}_{5}, \mathcal{S}_{5}\right)\right)
\end{gather*}
$$

such that

$$
\begin{equation*}
E(u)=u \text { in } X \tag{6.3}
\end{equation*}
$$

Before constructing the extension maps $E_{1}$ and $E_{2}$ we need to present a result on the commutation of blow-ups.

In this section we will modify the notation introduced in chapter 3 and apply the blow-ups symmetrically as in the construction of $\widetilde{X}_{5}$ (see Fig. 3.2 and Fig. 6.1):

$$
\widetilde{X}_{5} \rightarrow \tilde{X}_{3} \rightarrow \widetilde{X}_{2} \rightarrow \tilde{X}_{1} \rightarrow \tilde{X}
$$

Thus, $\widetilde{X}_{2}$ is a manifold with corners defined by the $1-2-3$ blow-up of $\Gamma$ followed by the $2-1-1$ blow-up of $\beta_{1}^{*} D \cap \partial \widetilde{X}_{1}$ with $\beta_{1} \circ \beta_{12}: \widetilde{X}_{2} \longrightarrow \widetilde{X}_{1}$ a corresponding blow-down map. Recall also that $\widetilde{X}_{3}$ is a manifold with corners obtained from $\widetilde{X}_{2}$ by the $2-1-0$ blow-up of $D^{(2)}=\left(\beta_{1} \circ \beta_{12}\right)^{*} D$ and $\beta_{23}: \widetilde{X}_{3} \longrightarrow \widetilde{X}_{2}$ is the corresponding blow-down map.

The problem of constructing the extension using $\beta_{1} \circ \beta_{12} \circ \beta_{23}$ is that the boundary of $X$ does not lift to a smooth hypersurface under this map. Thus we will need an alternative hierarchy of blow-ups.


Figure 6.1. The auxiliary blow-ups: $\tilde{X}_{a} \xrightarrow{\beta_{r a}} \widetilde{X}_{r} \xrightarrow{\beta_{2 r}} \widetilde{X}_{2}$

Let $D_{0}^{(2)}=D^{(2)} \cap \partial \widetilde{X}_{2}$ and $\left(x_{1}, x_{2}, x_{3}, y\right)$ be local coordinates given in $i$ ) of Proposition 3.3. Let $\left(z^{\prime}, y^{\prime}, s\right)$ be the corresponding projective coordinates near $D_{0}^{(2)}$ chosen so that $\beta_{12}\left(z^{\prime}, y^{\prime}, s\right) \mapsto\left(s^{2} z^{\prime}, s y^{\prime}, s\right)$. Then

$$
D^{(2)}=\left\{y^{\prime}=z^{\prime}=0\right\}, \quad D_{0}^{(2)}=\left\{s=y^{\prime}=z^{\prime}=0\right\}
$$

Let

$$
\mathbb{S}_{6-3-1_{ \pm}}^{2}=\left\{\omega \in \mathbb{R}^{3}: \omega_{1}^{2}+\omega_{2}^{4}+\omega_{3}^{12}=1, \quad \pm \omega_{3} \geq 0\right\}
$$

and define the space obtained by blowing up the two components of $D_{0}^{(2)}$ :

$$
\tilde{X}_{r}=\tilde{X}_{2} \backslash D_{0}^{(2)} \sqcup\left(S_{6-3-1_{+}}^{2} \times \mathbb{R}^{n-3}\right) \sqcup\left(S_{6-3-1_{-}}^{2} \times \mathbb{R}^{n-3}\right)
$$

with the $C^{\infty}$ structure obtained as before. The blow-down map near $D_{0}^{(2)}$ now takes form:

$$
\begin{gathered}
\beta_{2 r}: \tilde{X}_{r} \longrightarrow \tilde{X}_{2} \\
\beta_{2 r}(\omega, \rho)=\left(\rho^{6} \omega_{1}, \rho^{2} \omega_{2}, \rho \omega_{3}\right)=\left(z^{\prime}, y^{\prime}, s\right)
\end{gathered}
$$

If $D^{(3)}=\beta_{2 r}^{*} D^{(2)}$, we define the auxiliary manifold with corners $\widetilde{X}_{a}$ as the one obtained from $\widetilde{X}_{r}$ by blowing-up the submanifold $D^{(3)}$ with homogeneity $2-1-0$ and let $\beta_{r a}$ denote the corresponding blow-down map. Let

$$
\begin{gathered}
\beta_{2 a}: \tilde{X}_{a} \longrightarrow \tilde{X}_{2} \\
\beta_{2 a}=\beta_{2 r} \circ \beta_{r a}
\end{gathered}
$$

be the corresponding blow-down map.
Lemma 6.2. - If $Q=\left\{x_{2}^{2}-x_{1} x_{3}=0\right\}$ and $H=\left\{x_{1}=0\right\}$, then the smooth hypersurfaces $Q^{(2)}=\beta_{2}^{*} Q$ and $H^{(2)}=\beta_{2}^{*} H$ lift under $\beta_{2 a}$ to disjoint smooth hypersurfaces intersecting $\partial \widetilde{X}_{a}$ transversally.

Proof. - In the projective coordinates $\left(s, y^{\prime}, z^{\prime}\right)$ defined in a neighbourhood of $D^{(2)}$

$$
\begin{equation*}
Q^{(2)}=\left\{z^{\prime}=y^{\prime 2}\right\}, \quad H^{(2)}=\left\{z^{\prime}=0\right\} . \tag{6.4}
\end{equation*}
$$

As before we shall define projective coordinates for $\beta_{2 r}$. Thus let

$$
\begin{align*}
& \rho=s, \quad y^{\prime \prime}=\frac{y^{\prime}}{\rho^{3}}, \quad z^{\prime \prime}=\frac{z^{\prime}}{\rho^{6}} . \\
& \rho=\left|y^{\prime}\right|^{\frac{1}{3}}, \quad s^{\prime}=\frac{s}{\rho}, \quad z^{\prime \prime}=\frac{z^{\prime}}{\rho^{6}} .  \tag{6.5}\\
& \rho=\left|z^{\prime}\right|^{\frac{1}{6}}, \quad y^{\prime \prime}=\frac{y^{\prime}}{\rho^{3}}, \quad s^{\prime}=\frac{s}{\rho} .
\end{align*}
$$

Each valid in the region $\left|\omega_{i}\right|>0, i=1,2,3$. In each of these coordinate systems we obtain

$$
\begin{gathered}
\beta_{2 r}^{*} Q^{(2)}=\left\{z^{\prime \prime}=y^{\prime \prime 2}\right\}, \quad \beta_{2 r}^{*} H^{(2)}=\left\{z^{\prime \prime}=0\right\} \\
\beta_{2 r}^{*} Q^{(2)}=\left\{z^{\prime \prime}=1\right\}, \quad \beta_{2 r}^{*} H^{(2)}=\left\{z^{\prime \prime}=0\right\} \\
\beta_{2 r}^{*} Q^{(2)}=\left\{y^{\prime \prime}=1\right\}
\end{gathered}
$$

Hence in the first coordinate system $\beta_{2 r}^{*} Q^{(2)}$ and $\beta_{2 r}^{*} H^{(2)}$ are simply tangent and thus lift to smooth hypersurfaces in $\widetilde{X}_{a}$ by the $2-1-0$ blow-up. In the second and third coordinate systems $\beta_{2 r}^{*} Q^{(2)}$ and $\beta_{2 r}^{*} H^{(2)}$ are clearly smooth. This concludes the proof of the lemma.

Proposition 6.3. - Let $\left.{ }^{\circ}\right)_{a}=\varnothing\left(\beta_{2 a}^{*} Q^{(2)} \sqcup \beta_{2 a}^{*} H^{(2)} \sqcup \partial \widetilde{X}_{a}\right)$ be the Lie algebra of smooth vector fields in $\widetilde{X}_{a}$ tangent to $\beta_{2 a}^{*} Q^{(2)}, \beta_{2 a}^{*} H^{(2)}$ and $\partial \widetilde{X}_{a}$ and let $\left.{ }^{\circ}\right)_{3}=$ ©) $\left(\beta_{23}^{*} Q^{(2)} \sqcup \beta_{23}^{*} H^{(2)} \sqcup \partial \widetilde{X}_{3}\right)$ be the corresponding Lie algebra in $\widetilde{X}_{3}$. Let $\nu_{a}$ be the lift of the Lebesgue measure in $\mathbb{R}^{n}$ to $\widetilde{X}_{a}$. Then

$$
\begin{equation*}
\left(\beta_{23}\right)_{*} I_{k} L_{\nu_{3}}^{2}\left(\tilde{X}_{3}, \vartheta_{3}\right)=\left(\beta_{2 a}\right)_{*} I_{k} L_{\nu_{a}}^{2}\left(\tilde{X}_{a}, \vartheta_{a}\right) \tag{6.6}
\end{equation*}
$$

Proof. - Observe that in coordinates $\left(s, y^{\prime}, z^{\prime}\right)$ such that (6.4) holds, the map $F\left(s, y^{\prime}, z^{\prime}\right)=\left(s, y^{\prime}, z^{\prime}-{y^{\prime}}^{2}\right)$ is a diffeomorphism that preserves $D^{(2)}$ and maps $Q^{(2)}$ into $H^{(2)}$. The homogeneity of the variables $s, y^{\prime}$ and $z^{\prime}$ shows that $F$ lifts respectively under $\beta_{23}$ and $\beta_{2 a}$ to smooth diffeomorphisms mapping the lift of $Q^{(2)}$ into the corresponding lift of $H^{(2)}$. Since the pairs $\beta_{23}^{*} Q^{(2)}$ and $\beta_{23}^{*} H^{(2)}$ are disjoint, and so are $\beta_{a}^{*} Q^{(2)}$ and $\beta_{a}^{*} H^{(2)}$, we only need to prove (6.6) for the Lie algebras ${ }^{\circ} 0_{3}$ and $\left.{ }^{\circ}\right)_{a}$ replaced by $\left.{ }^{\circ}\right)_{3, H}=\vartheta\left(\beta_{23}^{*} H^{(2)} \sqcup \partial \widetilde{X}_{3}\right)$ and $\left.{ }^{\circ}\right)_{a, H}=\vartheta\left(\beta_{a}^{*} H^{(2)} \sqcup \partial \widetilde{X}_{a}\right)$.
The next step is to show that both sides of (6.6) have the same characterization in terms of singular vector fields.

Lemma 6.4. - Let $u \in L_{\nu_{2}}^{2}\left(\widetilde{X}_{2}\right)$. Then $u \in\left(\beta_{23}\right)_{*} I_{k} L_{\nu_{3}}^{2}\left(\widetilde{X}_{3}, \widetilde{\vartheta}_{3, H}\right)$ if and only if in the projective coordinates $\left(s, y^{\prime}, z^{\prime}\right)$ in which (6.4) holds

$$
\begin{equation*}
\left(s \partial_{s}, z^{\prime} \partial_{z^{\prime}},\left|z^{\prime}\right|^{\frac{1}{2}} \partial_{y^{\prime}}, y^{\prime} \partial_{y^{\prime}}\right)^{\alpha} u \in L_{\nu_{2}}^{2}\left(X_{2}\right), \quad|\alpha| \leq k \tag{6.7}
\end{equation*}
$$

Proof. - Indeed, if $u$ satisfies (6.7), then in projective coordinates

$$
\begin{equation*}
t=\left|y^{\prime}\right|, \quad z^{\prime \prime}=\frac{z^{\prime}}{t^{2}} \text { and } t=\left|z^{\prime}\right|^{\frac{1}{2}}, \quad y^{\prime \prime}=\frac{y^{\prime}}{t} \tag{6.8}
\end{equation*}
$$

the vector fields in (6.7) lift to

$$
\begin{equation*}
s \partial_{s}, z^{\prime \prime} \partial_{z^{\prime \prime}}, t \partial_{t},\left|z^{\prime \prime}\right|^{\frac{1}{2}}\left(s \partial_{s}-2 z^{\prime \prime} \partial_{z^{\prime \prime}}\right), \text { and } s \partial_{s}, t \partial_{t}, \partial_{y^{\prime \prime}} \tag{6.9}
\end{equation*}
$$

respectively, and these span ${ }^{\circ} \mathcal{O}_{3, H}$. Therefore $\left.\beta_{23}^{*} u \in I_{k} L_{\nu_{3}}^{2}\left(\widetilde{X}_{3},{ }^{\circ}\right)_{3, H}\right)$.
Conversely if $\left.\beta_{23}^{*} u \in I_{k} L_{\nu_{3}}^{2}\left(\tilde{X}_{3},{ }^{( }\right)_{3, H}\right)$, then $\beta_{23}^{*} u$ is stable under the application of the vector fields in (6.9) and therefore $u$ satisfies (6.7).

Similarly we obtain
Lemma 6.5. - Let $u \in L_{\nu_{2}}^{2}\left(\widetilde{X}_{2}\right)$. Then $\left.u \in\left(\beta_{2 a}\right)_{*} I_{k} L_{\nu_{a}}^{2}\left(\widetilde{X}_{a},{ }^{9}{ }^{( }\right) a, H\right)$ if and only if $u$ satisfies (6.7).
The two lemmas complete the proof of Proposition 6.3
We also need a corollary of Seeley's extension theorem, with the proof being immediate from the arguments of [46]. Let $\mathbb{R}_{+}^{n}=\left\{\left(x^{\prime}, x_{n}\right): x_{n}>0\right\}$ and for $s \in \mathbb{R}$, let $\bar{H}_{(s)}\left(\mathbb{R}_{+}^{n}\right)$ be the space of restrictions to $\mathbb{R}_{+}^{n}$ of elements in $H_{(s)}\left(\mathbb{R}^{n}\right)$

Proposition 6.6. - There exists a linear and continuous map

$$
S: L_{c}^{2}\left(\mathbb{R}_{+}^{n}\right) \longrightarrow L_{c}^{2}\left(\mathbb{R}^{n}\right)
$$

such that if $u \in L_{c}^{2}\left(\mathbb{R}_{+}^{n}\right)$ satisfies

$$
\begin{equation*}
Q_{1}\left(x^{\prime}, D_{x^{\prime}}\right) \ldots Q_{l}\left(x^{\prime}, D_{x^{\prime}}\right)\left(D_{x_{n}}^{r}\right)^{j} u \in \bar{H}_{(-(r-1) j)}\left(\mathbb{R}_{+}^{n}\right) \tag{6.10}
\end{equation*}
$$

for $Q_{i}\left(x^{\prime}, D_{x^{\prime}}\right) \in \operatorname{Diff}^{\bullet}\left(\mathbb{R}^{n-1}\right), 1 \leq i \leq m, j+l \leq k$ and $r \in \mathbb{N}, r \geq 1$. Then

$$
\begin{equation*}
Q_{1}\left(x^{\prime}, D_{x^{\prime}}\right) \ldots Q_{l}\left(x^{\prime}, D_{x^{\prime}}\right)\left(D_{x_{n}}^{r}\right)^{j} S(u) \in H_{(-(r-1) j)}\left(\mathbb{R}^{n}\right) \tag{6.11}
\end{equation*}
$$

We can now start the construction of the extension map and it is convenient to introduce the following notation:

$$
X_{i}=\beta_{i}^{*} X, \quad \partial X_{i}=\beta_{i}^{*} \partial X
$$

The $X_{i}$ 's are manifolds with corners and in the definition of the extendible $\left(H_{(s)}^{b}\right)$ and partial $\left(H_{(s, m)}^{b}\right) b$-Sobolev spaces (see Appendix B) are taken with respect to the boundary face given by $\partial X_{i}$. The first half of the theorem is given by

Proposition 6.7. - There exists a linear and continuous map

$$
\begin{equation*}
E_{1}: I_{k}^{b} L_{c}^{2}\left(X, \Lambda_{S} \sqcup \Lambda_{H}\right)+I_{k}^{b} L_{c}^{2}\left(X, \Lambda_{S} ; \Lambda_{S} \cap \Lambda_{F}\right) \longrightarrow\left(\beta_{5}\right)_{*}\left(I_{k} L_{\nu_{5}}^{2}\left(\widetilde{X}_{5}, \mathscr{S}_{5}\right)\right) \tag{6.12}
\end{equation*}
$$

such that

$$
\begin{equation*}
E_{1}(u)=u \text { in } X \tag{6.13}
\end{equation*}
$$

Proof. - Let us define the following Lagrangian variety

$$
\mathfrak{\varrho}=\left\{\Lambda_{S} \sqcup \Lambda_{H}, \Lambda_{S} \cap \Lambda_{F}\right\},
$$

where the marking of the union is disjoint from the second component. It easily follows that if $u$ satisfies the hypothesis of Proposition 6.7, then $u \in I_{k}^{b} L^{2}(X, \mathscr{Q})$. We shall use this to construct $\tilde{u}$ satisfying (6.13). Let us assume that $\operatorname{supp}(u)$ is contained in a small neighbourhood of $\Gamma$.

Proposition 3.3 guarantees the existence of smooth coordinates $(x, y)$ in a neighbourhood of $\Gamma$ such that

$$
\begin{equation*}
\widetilde{S}=\left\{x_{2}^{2}=x_{1} x_{3}\right\}, \quad B=\widetilde{F} \cap \widetilde{S}=\left\{x_{2}=x_{3}=0\right\}, \quad H=\left\{x_{1}=0\right\} \tag{6.14}
\end{equation*}
$$

and the defining function $\rho$ of $\partial X$, i.e $X=\{\rho>0\}, \partial X=\{\rho=0\}$ satisfies

$$
\begin{equation*}
\rho(x, y)=x_{2}+c_{1} x_{1}^{2}+c_{2} x_{2}^{2}+c_{3} x_{3}^{2}+\sum_{i<j} a_{i j} x_{i} x_{j}+O\left(|x|^{3}\right), \quad c_{1}, c_{3}>0 \tag{6.15}
\end{equation*}
$$

see Remark 3.4.
Direct computations show that $\Psi_{b}^{1, \bullet}(X, \mathscr{\varrho})$ is the $\Psi_{b}^{0, \bullet}(X)$ span of

$$
\begin{gather*}
V_{1}=3 x_{3} \partial_{x_{3}}+2 x_{2} \partial_{x_{2}}+x_{1} \partial_{x_{1}}, \quad V_{2}=2 x_{1} \partial_{x_{1}}+x_{2} \partial_{x_{2}} \\
V_{3}=\left(x_{2}^{2}-x_{1} x_{3}\right) \partial_{x_{1}}, \quad V_{4}=\left(x_{2}^{2}-x_{1} x_{3}\right) \partial_{x_{2}}  \tag{6.16}\\
V_{5}=\left(x_{2}^{2}-x_{1} x_{3}\right) \partial_{x_{3}}, \quad A L, L=\partial_{x_{2}}^{2}-4 \partial_{x_{1}} \partial_{x_{3}}
\end{gather*}
$$

where $A \in \Psi_{b}^{-1}(X)$ is elliptic. By Proposition 4.5, if $V=\left(V_{1}, \ldots, V_{5}\right)$, we find that

$$
\begin{equation*}
V^{\alpha} L^{a} u=\sum_{|\beta| \leq a}\left(\partial_{x_{1}}, \partial_{x_{2}}, \partial_{x_{3}}\right)^{\beta} u_{\beta}, \quad u_{\beta} \in L^{2}(X), \quad|\alpha|+a \leq k . \tag{6.17}
\end{equation*}
$$

The first step is to analyze the lift of (6.17) by $\beta_{1}$. Let $X_{1}=\left\{r^{-2} \beta_{1}^{*} \rho>0\right\}$, $V_{i}^{\prime}=\beta_{1}^{*} V_{i}, 1 \leq i \leq 5, L^{\prime}=r^{4} \beta_{1}^{*} L, W_{j}=r^{j} \beta_{1}^{*} \partial_{x_{j}}, 1 \leq j \leq 3$. We deduce from (6.17) that if $u_{1}=\beta^{*} u, V^{\prime}=\left(V_{1}^{\prime}, \ldots, V_{5}^{\prime}\right)$ then

$$
\begin{gather*}
V^{\prime \alpha} L^{\prime a} u_{1}=r^{4 a} \sum_{|\beta| \leq a} r^{-3 \beta_{1}-2 \beta_{2}-\beta_{3}} W^{\beta} u_{\beta}^{1}  \tag{6.18}\\
|\alpha|+a \leq k, \quad u_{\beta}^{1} \in r^{-3} L_{b}^{2}\left(X_{1}\right)
\end{gather*}
$$

Since $u$ is supported near $\Gamma$ we may assume that $r<1$ on $\operatorname{supp}\left(u_{1}\right)$. As $4 a-3 \beta_{1}-2 \beta_{2}-\beta_{3} \geq 0$, we obtain from (6.18) that

$$
\begin{equation*}
V^{\prime \alpha} L^{\prime a} u_{1} \in r^{-3} \bar{H}_{(-a)}^{b}\left(X_{1}\right) \tag{6.19}
\end{equation*}
$$

where $L_{b}^{2}$ and $H_{(s)}^{b}$ are the $b$-Sobolev spaces defined in Appendix B. Then we analyze (6.19) in projective coordinates

$$
\begin{align*}
& r=\left|x_{3}\right|^{\frac{1}{3}}, \quad Y=\frac{x_{2}}{r^{2}}, \quad Z=\frac{x_{1}}{r} \\
& r=\left|x_{2}\right|^{\frac{1}{2}}, \quad X=\frac{x_{3}}{r^{3}}, \quad Z=\frac{x_{1}}{r}  \tag{6.20}\\
& r=\left|x_{1}\right|, \quad X=\frac{x_{3}}{r^{3}}, \quad Y=\frac{x_{2}}{r^{2}} .
\end{align*}
$$

In the third set of coordinates we obtain from (6.19)

$$
\begin{equation*}
\left(r \partial_{r}, 2 X \partial_{X}+Y \partial_{Y}\right)^{\alpha}\left[\left(\partial_{Y}+2 Y \partial_{X}\right) \partial_{Y}\right]^{a} u_{1} \in r^{-3} \bar{H}_{(-a)}^{b}\left(X_{1}\right), \quad|\alpha|+a \leq k \tag{6.21}
\end{equation*}
$$

and

$$
\begin{equation*}
\rho_{1}=r^{-2} \beta_{1}^{*} \rho=Y+c_{1}+O(r) \tag{6.22}
\end{equation*}
$$

Thus for small $r, Y \neq 0$ near $\partial X_{1}$. The operators in (6.21) span the space of totally characteristic operators in $\Psi_{b}^{1}\left(\widetilde{X}_{1},{ }^{b} £\right)$ where ${ }^{b}{ }_{\varrho}$ is the marked Lagrangian variety formed by ${ }^{b} N^{*} \widetilde{S}^{(1)} \sqcup^{b} N^{*} H^{(1)}$ marked by ${ }^{b} N^{*} \widetilde{S^{(1)}} \cap^{b} N^{*} B^{(1)}$.
Near $\partial X_{1}, \rho=Y, x^{\prime}=X / Y^{2}$, give a smooth coordinate system in which (6.21) can be written as

$$
\begin{gather*}
\left(r \partial_{r}, \partial_{\rho}\right)^{\alpha}\left[x^{\prime}\left(1-x^{\prime}\right) \partial_{x^{\prime}}^{2}\right]^{a} u_{1} \in r^{-3} \bar{H}_{(-a)}^{b}\left(X_{1}\right) \\
\text { and }  \tag{6.23}\\
\partial X_{1}=\left\{\rho+c_{1}+O(r)=0\right\}
\end{gather*}
$$

If $\rho^{\prime}=0$ defines the boundary then $\rho^{\prime}=\rho+c_{1}+O(r)$ and we can use it as a new coordinate. Then we deduce from (6.23) that

$$
\begin{equation*}
\left(r \partial_{r}, \partial_{\rho^{\prime}}\right)^{\alpha}\left(x^{\prime}\left(1-x^{\prime}\right) \partial_{x^{\prime}}^{2}\right)^{a} u_{1} \in r^{-3} \bar{H}_{(-a)}^{b}\left(X_{1}\right), \quad \partial X_{1}=\left\{\rho^{\prime}=0\right\} \tag{6.24}
\end{equation*}
$$

Since $\partial_{\rho^{\prime}}$ is transversal to $\partial X_{1}$ we deduce from Proposition 6.6 , or rather its easy modification to the case of a manifold with corners, that there exists $\tilde{u}$ such that

$$
\begin{equation*}
\tilde{u}=u \text { in } X_{1}, \quad\left(r \partial_{r}, \partial_{\rho^{\prime}}\right)^{\alpha}\left(x^{\prime}\left(1-x^{\prime}\right) \partial_{x^{\prime}}^{2}\right)^{a} \tilde{u} \in r^{-3} \bar{H}_{(-a)}^{b}\left(\tilde{X}_{1}\right) \tag{6.25}
\end{equation*}
$$

Now we proceed by a microlocal partition of unity of ${ }^{b} T^{*} \widetilde{X}_{1}$. In the region where either $r \partial_{r}$ or $\partial_{\rho^{\prime}}$ is elliptic it follows from (6.25) that $\tilde{u}$ is conormal to $\partial \widetilde{X}_{1}$. In the region where $\partial_{x^{\prime}}$ is elliptic we obtain from (6.25) that

$$
\begin{equation*}
\left(r \partial_{r}, \partial_{\rho^{\prime}}\right)^{\alpha}\left(x^{\prime}\left(1-x^{\prime}\right) \partial_{x^{\prime}}\right)^{a} \tilde{u} \in r^{-3} L_{b}^{2}\left(\widetilde{X}_{1}\right) \tag{6.26}
\end{equation*}
$$

Going back to the original coordinates this provides the condition required in the extension.

In the region where the coordinates $(r, X, Z)$ can be used

$$
\partial X_{1}=\left\{1+c_{1} Z^{2}+O(r)=0\right\}
$$

Thus for small $r, Z \neq 0$ near $\partial X_{1}$ and therefore the set of projective coordinates $(r, Y, X)$ can also be used there and the extension is constructed as above.
In the first set of projective coordinates in (6.20) we obtain,

$$
\begin{gather*}
\left(r \partial_{r}, 2 Z \partial_{Z}+Y \partial_{Y}, r^{2}\left(Z-Y^{2}\right) \partial_{Y}, r^{3}\left(Z-Y^{2}\right) \partial_{Z}\right)^{\alpha} \\
{\left[\left(\partial_{Y}+2 Y \partial_{Z}\right) \partial_{Y}\right]^{a} u_{1} \in r^{-3} \bar{H}_{(-a)}^{b}\left(X_{1}\right) .} \tag{6.27}
\end{gather*}
$$

and the lift of the boundary of $\partial X$ is given by

$$
\begin{gathered}
\partial X_{1}=\left\{Y+c_{1} Z^{2}+a_{12} r Y Z+a_{13} r^{2} Z+b_{1} r Z^{3}+b_{2} r^{2} Y Z^{2}+f(Y, Z, r) r^{3}=0\right\} \\
f \in C^{\infty}, \quad f(0,0,0)=0
\end{gathered}
$$

Since we are only concerned with the region $r \sim 0$, we find that if $|Y|>0$ near $\partial X_{1}$, then $|Z|>0$ near $\partial X_{1}$. Therefore in this case the extension map is constructed in the third set of coordinates. Thus we may restrict our analysis to a small neighbourhood of $Y=Z=0$.
To start we need to analyze the lift of (6.27) under the $2-1-1$ blow-down map. For that we use projective coordinates

$$
\begin{align*}
& s=r, \quad y^{\prime}=\frac{Y}{s}, \quad z^{\prime}=\frac{Z}{s^{2}} \\
& s=|Y|, \quad r^{\prime}=\frac{r}{s}, \quad z^{\prime}=\frac{Z}{s^{2}}  \tag{6.28}\\
& s=|Z|^{\frac{1}{2}}, \quad r^{\prime}=\frac{r}{s}, \quad y^{\prime}=\frac{Y}{s}
\end{align*}
$$

In the third set of coordinates of (6.28)

$$
\begin{equation*}
\partial X_{2}=\left\{y^{\prime}+O\left(s^{3}\right)=0\right\} \tag{6.29}
\end{equation*}
$$

We deduce from (6.27) that $u_{2}=\beta_{12}^{*} u_{1}$ satisfies

$$
\begin{equation*}
\left(r \partial_{r^{\prime}}, s \partial_{s}\right)^{\alpha}\left[\left(1 \mp y^{\prime 2}\right) \partial_{y^{\prime}}^{2}\right]^{a} u_{2} \in{r^{\prime}}^{-3} s^{-\frac{9}{2}} \bar{H}_{(-a)}^{b}\left(X_{2}\right), \quad|\alpha|+a \leq k \tag{6.30}
\end{equation*}
$$

Since $r^{\prime}$ and $s$ are small, $y^{\prime} \sim 0$ near $\partial X_{2}$ and thus $\left|1 \mp{y^{\prime}}^{2}\right|>0$ there. Therefore near $\partial X_{2}$

$$
\begin{equation*}
\left(r^{\prime} \partial_{r^{\prime}}, s \partial_{s}\right)^{\alpha} \partial_{y^{\prime}}^{2 a} u_{2} \in r^{\prime-3} s^{-\frac{9}{2}} \bar{H}_{(-a)}^{b}\left(X_{2}\right), \quad|\alpha|+a \leq k \tag{6.31}
\end{equation*}
$$

Let $y^{\prime \prime}=y^{\prime}+O\left(s^{3}\right)$ be a defining function of $\partial X_{2}$. We deduce from (6.31) that

$$
\begin{gather*}
\left(r^{\prime} \partial_{r^{\prime}}+f_{1} \partial_{y^{\prime \prime}}, s \partial_{s}+f_{2} \partial_{y^{\prime \prime}}\right)^{\alpha} \partial_{y^{\prime \prime}}^{2 a} u_{2} \in r^{\prime-3} s^{-\frac{9}{2}} \bar{H}_{(-a)}^{b}\left(X_{2}\right), \quad|\alpha|+a \leq k  \tag{6.32}\\
\partial X_{2}=\left\{y^{\prime \prime}=0\right\}
\end{gather*}
$$

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where $f_{1}$ and $f_{2}$ are smooth functions. Thus, completing the squares, we see from (6.32) that

$$
\begin{equation*}
\left(\left(r^{\prime} \partial_{r^{\prime}}\right)^{2},\left(s \partial_{s}\right)^{2}, \partial_{y^{\prime \prime}}^{2}\right)^{\alpha} u_{2} \in r^{\prime-3} s^{-\frac{9}{2}} \bar{H}_{(-|\alpha|)}^{b}\left(X_{2}\right), \quad|\alpha| \leq k \tag{6.33}
\end{equation*}
$$

It follows from Proposition B.1, applied to the noncharacteristic operator $\partial_{y^{\prime \prime}}^{2 a}$, that

$$
\begin{equation*}
\left(\left(r^{\prime} \partial_{r^{\prime}}\right)^{2},\left(s \partial_{s}\right)^{2}\right)^{\alpha} \partial_{y^{\prime \prime}}^{a} u_{2} \in r^{\prime-3} s^{-\frac{9}{2}} \bar{H}_{(-|\alpha|)}^{b}\left(X_{2}\right), \quad|\alpha|+a \leq k \tag{6.34}
\end{equation*}
$$

Proposition 6.6 gives $\tilde{u}_{2} \in{r^{\prime}}^{-3} s^{-\frac{9}{2}} L_{b}^{2}\left(\tilde{X}_{2}\right)$ such that

$$
\begin{equation*}
\left.\left(\left(r^{\prime} \partial_{r^{\prime}}\right)^{2},\left(s \partial_{s}\right)^{2}\right)\right)^{\alpha} \partial_{y^{\prime \prime}}^{a} u_{2} \in{r^{\prime}}^{-3} s^{-\frac{9}{2}} H_{(-|\alpha|)}^{b}\left(\widetilde{X}_{2}\right), \quad|\alpha|+a \leq k \tag{6.35}
\end{equation*}
$$

Since at any point $q \in{ }^{b} T^{*} \widetilde{X}_{2}$ one of the three operators $s \partial_{s}, r^{\prime} \partial_{r^{\prime}}$ or $\partial_{y^{\prime \prime}}$ is elliptic one deduces from (6.35) that

$$
\widetilde{u}_{2} \in r^{\prime-3} s^{-\frac{9}{2}} H_{(k)}^{b}\left(\widetilde{X}_{2}\right)
$$

Simple computations show that, for $r^{\prime}$ and $s^{\prime}$ small, $\partial X_{2}$ does not intersect the region where the second coordinate system in (6.28) holds.
Next we consider the region where the first coordinate system in (6.28) is used. We find that in these coordinates

$$
\begin{equation*}
\partial X_{2}=\left\{y^{\prime}+c_{1} s^{3}+s^{3} f\left(r^{\prime}, z^{\prime}, s\right)=0\right\}, \quad f \in C^{\infty}, \quad f(0,0,0)=0 \tag{6.36}
\end{equation*}
$$

and we obtain from (6.19) that

$$
\begin{equation*}
\left(s \partial_{s}, 2 z^{\prime} \partial_{z^{\prime}}+y^{\prime} \partial_{y^{\prime}}\right)^{\alpha}\left[\left(\partial_{y^{\prime}}+2 y^{\prime} \partial_{z^{\prime}}\right) \partial_{y^{\prime}}\right]^{m} u_{2} \in s^{-\frac{9}{2}} \bar{H}_{(-m)}^{b}\left(X_{2}\right), \quad|\alpha|+m \leq k \tag{6.37}
\end{equation*}
$$

To analyze $\widetilde{X}_{3}$ we need to blow-up $D^{(2)}$ with homogeneity $1-2$. However as we already mentioned $\partial X_{2}$ lifts under $\beta_{23}$ to a singular hypersurface. This is where we use the result on the commutation of blow-ups proved in Proposition 6.3. We shall prove that there exists an extension of $\beta_{2 a}^{*} u_{2}$ into the conormal space $I_{k} L_{\nu_{a}}^{2}\left(X_{a},{ }^{,}{ }_{a}\right)$. It follows from Proposition 6.6 that this gives an extension of $u_{2}$ to the conormal space.
In the first set of coordinates (6.5) we obtain that $u_{3}=\beta_{2 r}^{*} u_{2}$ satisfies

$$
\begin{gather*}
\left(\rho \partial_{\rho}, 2 z^{\prime \prime} \partial_{z^{\prime \prime}}+y^{\prime \prime} \partial_{y^{\prime \prime}}\right)^{\alpha}\left[\left(\partial_{y^{\prime \prime}}+2 y^{\prime \prime} \partial_{z^{\prime \prime}}\right) \partial_{y^{\prime \prime}}\right]^{m} u_{3} \in \rho^{-9} \bar{H}_{(-m)}^{b}\left(X_{r}\right), \quad|\alpha|+m \leq k . \\
\quad \text { and } \\
\partial X_{r}=\left\{y^{\prime \prime}+c_{1}+O(\rho)=0\right\} . \tag{6.38}
\end{gather*}
$$

Thus for small $\rho, y^{\prime \prime} \neq 0$ near $\partial X_{r}$. Then we blow-up $D^{(r)}=\beta_{2 r}^{*} D^{(2)}=\left\{y^{\prime \prime}=z^{\prime \prime}=\right.$ $0\}$ with homogeneity $2-1-0$. In projective coordinates

$$
\begin{equation*}
t=\left|y^{\prime \prime}\right|, \quad \tilde{z}=\frac{z^{\prime \prime}}{t^{2}} \text { and } \tilde{y}=\frac{y^{\prime \prime}}{t}, \quad t=\left|z^{\prime \prime}\right|^{\frac{1}{2}} \tag{6.39}
\end{equation*}
$$

The condition (6.37) gives, for $|\alpha|+m \leq k$

$$
\begin{align*}
& \left(\rho \partial_{\rho}, t \partial_{t}\right)^{\alpha}\left[\tilde{z}(\tilde{z}-1) \partial_{\tilde{z}}^{2}\right]^{m} \beta_{r a}^{*} u_{3} \in t^{-\frac{3}{2}} \rho^{-9} \bar{H}_{(-m)}\left(X_{a}\right) \\
& \left(\rho \partial_{\rho}, t \partial_{t}\right)^{\alpha}\left[\tilde{y}(\tilde{y}-1) \partial_{\tilde{y}}^{2}\right]^{m} \beta_{r a}^{*} u_{3} \in t^{-\frac{3}{2}} \rho^{-9} \bar{H}_{(-m)}\left(X_{a}\right) \tag{6.40}
\end{align*}
$$

Hence, away from $\partial X_{a}$, we conclude by a simple microlocal partition of unity argument that $\beta_{r a}^{*} u_{3} \in I_{k} L_{\nu_{a}}\left(X_{a},{ }^{\left.()_{a}\right)}\right.$ ).
Since $y^{\prime \prime} \neq 0$ near $\partial X_{r}$, the first set of coordinates in (6.39) can be used and

$$
\partial X_{a}=\left\{t+c_{1}+O(\rho)=0\right\}
$$

We deduce from (6.40) that

$$
\begin{equation*}
\left(\rho \partial_{\rho}, \partial_{t}\right)^{\alpha}\left[\tilde{z}(\tilde{z}-1) \partial_{\tilde{z}}^{2}\right]^{m} \beta_{r a}^{*} u_{3} \in \rho^{-9} \bar{H}_{(-m)}\left(X_{a}\right) \tag{6.41}
\end{equation*}
$$

Since $\partial_{t}$ is transversal to $\partial X_{3}$, one deduces from Proposition 6.6 that there exists an extension $\tilde{u}_{3}$ of $\beta_{r a}^{*} u_{3}$ to $\widetilde{X}_{a}$ such that

$$
\begin{equation*}
\left(\rho \partial_{\rho}, t \partial_{t}\right)^{\alpha}\left[\tilde{z}(\tilde{z}-1) \partial_{\tilde{z}}^{2}\right]^{m} \tilde{u}_{3} \in \rho^{-\frac{13}{2}} \bar{H}_{(-m)}\left(\tilde{X}_{a}\right) \tag{6.42}
\end{equation*}
$$

Next we proceed by a microlocal partition of unity. In the region where either $\rho \partial_{\rho}$ or $t \partial_{t}$ is elliptic $\tilde{u}_{3}$ is conormal to the boundary of $\tilde{X}_{a}$. In the region where $\partial_{\tilde{z}}$ is elliptic we find that

$$
\begin{equation*}
\left(\rho \partial_{\rho}, t \partial_{t}\right)^{\alpha}\left[\tilde{z}(\tilde{z}-1) \partial_{\tilde{z}}\right]^{a} \tilde{u}_{3} \in \rho^{-9} L_{b}^{2}\left(\tilde{X}_{a}\right) \tag{6.43}
\end{equation*}
$$

This shows that $\tilde{u}_{3}$ is conormal to the lifting of the cone and the plane.
Simple computations show that for $\rho$ and $s^{\prime}$ small, the boundary $\partial X_{r}$ does not intersect the region where the second coordinate system in (6.5) holds. In the third set of coordinates

$$
\begin{gather*}
\left(\rho \partial_{\rho}, s \partial_{s}\right)^{\alpha}\left[\left(1 \mp y^{\prime \prime 2}\right) \partial_{y^{\prime \prime}}^{2}\right]^{m} u_{3} \in s^{\prime-\frac{9}{2}} \rho^{-9} \bar{H}_{(-m)}^{b}\left(X_{r}\right), \quad m+|\alpha| \leq k \\
\quad \text { and }  \tag{6.44}\\
\partial X_{r}=\left\{y^{\prime \prime}+O(\rho)=0\right\}
\end{gather*}
$$

Away from $\partial X_{r}$, a microlocal partition of unity argument gives that $u_{3} \in$ $I_{k} L_{\nu_{r}}\left(X_{r},{ }^{\circ} \vartheta_{r}\right)$. On the other hand $1 \mp y^{\prime \prime} \neq 0$ near $\partial X_{r}$. Thus it follows from Proposition 6.6 that there exists an extension $\tilde{u}_{3}$ of $u_{3}$ to $\widetilde{X}_{r}$ such that

$$
\begin{equation*}
\left(\rho \partial_{\rho}, s \partial_{s}\right)^{\alpha} \partial_{y^{\prime \prime}}^{2 m} \tilde{u}_{3} \in s^{\prime-\frac{9}{2}} \rho^{-9} \bar{H}_{(-m)}^{b}\left(X_{r}\right), \quad m+|\alpha| \leq k \tag{6.45}
\end{equation*}
$$

Now a simple microlocal partition of unity argument analogous to the one used above shows that $\tilde{u}_{3} \in \rho^{-11} s^{\prime-\frac{5}{2}} H_{(k)}^{b}\left(X_{r}\right)$. This concludes the proof of Proposition 6.7.
Let us observe that the linearity and continuity of $E_{1}$ follow from its construction and the linearity and continuity of the Seeley map $S$ in Proposition 6.6.

Next we consider the Lagrangians corresponding to the direct and reflected fronts, that is the second part of Theorem 6.1.

Proposition 6.8. - There exists a linear and continuous map

$$
\begin{equation*}
E_{2}: I_{k}^{b} L^{2}\left(X, \Lambda_{F} \sqcup \Lambda_{R}\right) \longrightarrow\left(\beta_{5}\right)_{*}\left(I_{k} L_{\nu_{5}}^{2}\left(\widetilde{X}_{5}, \mathscr{S}_{5}\right)\right) \tag{6.46}
\end{equation*}
$$

such that

$$
\begin{equation*}
E_{2}(u)=u \text { in } X \tag{6.47}
\end{equation*}
$$

Proof. - We can assume that $\operatorname{supp}(u)$ is contained in a small neighbourhood of $\Gamma$. Let $\left(x_{1}, x_{2}, x_{3}, y\right)$ be local coordinates in a neighbourhood of $\Gamma$ given by $i i$ ) of Proposition 3.3. Hence,

$$
\begin{equation*}
\widetilde{R}=\left\{x_{2}^{3}=x_{3}^{2}\right\}, \quad \widetilde{F}=\left\{2 x_{3}+x_{1}^{3}-3 x_{1} x_{2}=0\right\} \tag{6.48}
\end{equation*}
$$

and the boundary of $X$ is given by

$$
\begin{equation*}
\partial X=\left\{x_{2}=\frac{1}{4} x_{1}^{2}+g\right\}, \quad g \in M_{3}^{3-2-1} \tag{6.49}
\end{equation*}
$$

In these coordinates the space $\Psi_{b}^{1, \bullet}\left(X, \Lambda_{F} \sqcup \Lambda_{R}\right)$ is spanned over $\Psi_{b}^{0, \bullet}(X)$ by

$$
\begin{gather*}
A P, A L, P=4 \partial_{x_{2}}^{2}-9 x_{2} \partial_{x_{3}}^{2}-6 \partial_{x_{1}} \partial_{x_{3}}, \quad L=\left(2 \partial_{x_{2}}+3 x_{1} \partial_{x_{3}}\right) \partial_{x_{1}}  \tag{6.50}\\
V_{1}=3 x_{3} \partial_{x_{3}}+2 x_{2} \partial_{x_{2}}+x_{1} \partial_{x_{1}}, \quad \partial_{y_{j}}
\end{gather*}
$$

where $A \in \Psi_{b}^{-1}(X)$ is elliptic. To see that one needs to observe that in coordinates $\left(x_{1}, x_{2}, x_{3}, y, \xi_{1}, \xi_{2}, \xi_{3}, \eta\right)$ in $T^{*} \mathbb{R}^{n} \backslash 0$

$$
\begin{aligned}
& \Lambda_{R}=\left\{4 \xi_{2}^{2}-9 x_{2} \xi_{3}^{2}=0, \quad 3 x_{3} \xi_{3}+2 x_{2} \xi_{2}=0, \quad \xi_{1}=0, \quad \eta=0, \quad \xi_{3} \neq 0\right\} \\
\Lambda_{F}= & \left\{2 \xi_{2}+3 x_{1} \xi_{3}=0, \quad 2 \xi_{1}-3\left(x_{1}^{2}-x_{2}\right) \xi_{3}=0,3 x_{3} \xi_{3}+2 x_{2} \xi_{2}+x_{1} \xi_{1}=0, \quad \eta=0\right\}
\end{aligned}
$$

We then notice that
$\Lambda_{F} \cup \Lambda_{R} \subset M=\left\{4 \xi_{2}^{2}-9 x_{2} \xi_{3}^{2}-6 \xi_{1} \xi_{3}=0, \quad 3 x_{3} \xi_{3}+2 x_{2} \xi_{2}+x_{1} \xi_{1}=0, \eta=0, \xi_{3} \neq 0\right\}$
Then $M$ is a smooth submanifold of $T^{*} \mathbb{R}^{n} \backslash 0$ and in $M$

$$
\Lambda_{F}=\left\{\xi_{1}=0\right\}, \quad \Lambda_{R}=\left\{2 \xi_{2}+3 x_{1} \xi_{2}=0\right\}
$$

Therefore if $f \in C^{\infty}\left(T^{*} \mathbb{R}^{n} \backslash 0\right)$ vanishes on $\Lambda_{F} \cup \Lambda_{R}$

$$
\begin{gathered}
f=a_{1}\left(4 \xi_{2}^{2}-9 x_{2} \xi_{3}^{2}-6 \xi_{1} \xi_{3}\right)+a_{2}\left(3 x_{3} \xi_{3}+2 x_{2} \xi_{2}+x_{1} \xi_{1}\right)+a_{3}\left(2 \xi_{2}+3 x_{1} \xi_{2}\right) \xi_{1}+ \\
\sum a_{j} \eta_{j}, \quad a_{i} \in C^{\infty}\left(T^{*} \mathbb{R}^{n} \backslash 0\right)
\end{gathered}
$$

This shows that if $B \in \Psi_{b}^{m, k}\left(X, \Lambda_{F} \sqcup \Lambda_{R}\right)$ then $B$ is in the span of $P, L$ and $V_{1}$. Here and in what follows we should neglect the trivial generators $\partial_{y_{j}}$. It will also be useful for us to observe that

$$
V_{2}=3 x_{2}^{2} \partial_{x_{3}}+2 x_{3} \partial_{x_{2}}+\left(2 x_{2}-x_{1}^{2}\right) \partial_{x_{1}} \in \Psi_{b}^{1,1}\left(X, \Lambda_{F} \sqcup \Lambda_{R}\right) .
$$

Therefore if $u \in I_{k}^{b}\left(X, \Lambda_{F} \sqcup \Lambda_{R}\right)$ then Proposition 4.5 shows that

$$
\begin{equation*}
(P, L)^{a}\left(V_{1}, V_{2}\right)^{\alpha} u \in \bar{H}_{(-|a|)}(X), \quad|a|+|\alpha| \leq k \tag{6.51}
\end{equation*}
$$

where $V_{2}$ is used as a generator only for convenience.
To construct the extension map we need to examine the lift of (6.51) under the blow-down map $\beta_{1}$. Thus let

$$
\begin{align*}
P_{1} & =r^{4} \beta_{1}^{*} P, \quad L_{1}=r^{3} \beta_{1}^{*} L  \tag{6.52}\\
V_{i}^{\prime}=\beta_{1}^{*} V_{i}, i & =1,2, \quad W_{j}=r^{j} \beta_{1}^{*} \partial_{x_{j}}, j=1,2,3
\end{align*}
$$

We deduce from (6.51) that, for $u_{1}=\beta_{1}^{*} u$ and $W=\left(W_{1}, W_{2}, W_{3}\right)$

$$
\begin{gather*}
\left(P_{1}, L_{1}\right)^{a}\left(V_{1}^{\prime}, V_{2}^{\prime}\right)^{\alpha} u_{1}=r^{4 a_{1}+3 a_{2}} \sum_{|\beta| \leq|a|} r^{-3 \beta_{1}-2 \beta_{2}-\beta_{3}} W^{\beta} u_{\beta}  \tag{6.53}\\
u_{\beta} \in r^{-3} L_{b}^{2}\left(X_{1}\right), \quad|a|+|\alpha| \leq k
\end{gather*}
$$

Since $u$ is supported in a small neighbourhood of $\Gamma$, we may assume that $r<1$ on $\operatorname{supp}\left(u_{1}\right)$. Therefore,

$$
\begin{equation*}
\left(P_{1}, L_{1}\right)^{a}\left(V_{1}^{\prime}, V_{2}^{\prime}\right)^{\alpha} u_{1} \in r^{-3} \bar{H}_{(-|a|)}^{b}\left(X_{1}\right), \quad|\alpha|+|a| \leq k \tag{6.54}
\end{equation*}
$$

Since $x_{2} \geq 0$ on $\widetilde{R}$, we first consider the region where the projective coordinates $r=x_{2}^{\frac{1}{2}}, \quad X=x_{3} / r^{3}, \quad Z=x_{1} / r^{3}$ can be used. In terms of these we find that:

$$
\begin{gather*}
P_{1}=\left(r \partial_{r}-3 X \partial_{X}-Z \partial_{Z}\right)^{2}-9 \partial_{X}^{2}-6 \partial_{X} \partial_{Z}, \quad L_{1}=\left(r \partial_{r}+3(Z-X) \partial_{X}-Z \partial_{Z}\right) \partial_{Z} \\
V_{1}^{\prime}=r \partial_{r}, \quad V_{2}^{\prime}=r\left(3\left(1-X^{2}\right) \partial_{X}+\left(2-Z^{2}-X Z\right) \partial_{Z}\right) \tag{6.55}
\end{gather*}
$$

where we neglected the $r \partial_{r}$ terms in $V_{2}^{\prime}$. The boundary of $X_{1}$ and the lifted hypersurfaces are defined by

$$
\begin{gather*}
\partial X_{1}=\left\{Z^{2}=4+r \phi_{1}\right\}, \quad \phi_{1} \in C^{\infty}\left(\widetilde{X}_{1}\right) \\
\beta_{1}^{*} \widetilde{F}=\left\{2 X+Z^{3}-3 Z=0\right\}, \quad \beta_{1}^{*} \widetilde{R}=\left\{X^{2}-1=0\right\} \tag{6.56}
\end{gather*}
$$

One can write

$$
\beta_{1}^{*} \widetilde{F}=\left\{2(X \mp 1)+(Z \mp 1)^{2}(Z \mp 2)=0\right\}
$$



Figure 6.2. Intersections of $\beta^{*} \widetilde{F}$ and $\beta^{*} \widetilde{R}$
which shows that $\beta_{1}^{*} \widetilde{F}$ and $\beta_{1}^{*} \widetilde{R}$ are simply tangent along $B^{(1)}=\{X=Z= \pm 1\}$ and intersect transversally at $K^{(1)}=\{X= \pm 1, Z=\mp 2\}$, see Fig 6.2.
Since $r$ is small in $\operatorname{supp}\left(u_{1}\right)$ we may write

$$
\begin{equation*}
\partial X_{1}=\left\{Z=-2+r \phi_{2}\right\} \cup\left\{Z=2+r \phi_{3}\right\}, \quad \phi_{2}, \phi_{3} \in C^{\infty}\left(\widetilde{X}_{1}\right) . \tag{6.57}
\end{equation*}
$$

Let

$$
\begin{equation*}
P_{1}^{\prime}=\left(3 X \partial_{X}+Z \partial_{Z}\right)^{2}-9 \partial_{X}^{2}-6 \partial_{X} \partial_{Z}, \quad L_{1}^{\prime}=\left(3(Z-X) \partial_{X}-Z \partial_{Z}\right) \partial_{Z} \tag{6.58}
\end{equation*}
$$

Direct computations show that, modulo lower order terms,

$$
\begin{gather*}
P_{1}^{\prime}+Z L_{1}^{\prime}=Q_{1}^{\prime}=3 U_{2} \partial_{X} \\
U_{2}=\frac{1}{r} V_{2}^{\prime}=3\left(1-X^{2}\right) \partial_{X}+\left(2-Z^{2}-X Z\right) \partial_{Z} \tag{6.59}
\end{gather*}
$$

Hence we obtain from (6.55),(6.58) and (6.59) that

$$
\begin{equation*}
L_{1}^{a_{1}}\left(\partial_{X} U_{2}\right)^{a_{2}}\left(V_{1}^{\prime}, V_{2}^{\prime}\right)^{\alpha} u_{1} \in r^{-3} \bar{H}_{(-|a|)}^{b}\left(X_{1}\right), \quad|a|+|\alpha| \leq k \tag{6.60}
\end{equation*}
$$

We shall prove that in fact

$$
\begin{equation*}
L_{1}^{a_{1}} U_{2}^{a_{2}}\left(V_{1}^{\prime}, V_{2}^{\prime}\right)^{\alpha} u_{1} \in r^{-3} \bar{H}_{\left(-a_{1}\right)}^{b}\left(X_{1}\right), \quad|a|+|\alpha| \leq k \tag{6.61}
\end{equation*}
$$

First we observe that for $a_{1}+\alpha_{1} \leq k$

$$
\begin{equation*}
L_{1}^{a_{1}} U_{2}^{\alpha_{2}}\left(V_{1}^{\prime}\right)^{\alpha_{1}} u_{1} \in r^{-3} \bar{H}_{\left(-a_{1}-\alpha_{2}\right)}^{b}\left(X_{1}\right), \quad a_{1}+\alpha_{1} \leq k \tag{6.62}
\end{equation*}
$$

From (6.60) and the fact that $V_{2}^{\prime}=r U_{2}$, we deduce for $l+|\alpha| \leq k$

$$
\begin{align*}
& \partial_{X}^{\alpha_{2}} L_{1}^{l}\left(V_{1}^{\prime}, U_{2}\right)^{\alpha} u_{1}, \quad\left(r \partial_{r}\right)^{\alpha_{2}} L_{1}^{l}\left(V_{1}^{\prime}, U_{2}\right)^{\alpha} u_{1}, \\
& \left(r \partial_{Z}\right)^{\alpha_{2}} L_{1}^{l}\left(V_{1}^{\prime}, U_{2}\right)^{\alpha} u_{1} \in r^{-3} \bar{H}_{\left(-l-\alpha_{2}\right)}^{b}\left(X_{1}\right) \tag{6.63}
\end{align*}
$$

In the neighborhood where $Z \sim \pm 2$ we consider a change of variables

$$
\Psi(r, X, Z)=\left(r, X, Z-r \phi_{i}\right), \quad i=2,3
$$

$\Psi$ is a diffeomorphism preserving $\{r=0\}$ and $X_{1}^{\#}=\Psi\left(X_{1}\right)=\{Z \leq 2\}$. Let $U_{2}^{\#}=\Psi_{*} U_{2}, V_{1}^{\#}=\Psi_{*} V_{1}^{\prime}, L_{1}^{\#}=\Psi_{*} L_{1}$ and $u_{1}^{\#}=\Psi_{*} u_{1}$. We deduce from (6.63) that for $l+|\alpha| \leq k$

$$
\begin{equation*}
\partial_{X}^{\alpha_{2}} L_{1}^{\# l}\left(U_{2}^{\#}, V_{1}^{\#}\right)^{\alpha} u_{1}^{\#}, \quad\left(r \partial_{r}\right)^{\alpha_{2}} L_{1}^{\#^{l}}\left(U_{2}^{\#}, V_{1}^{\#}\right)^{\alpha} u_{1}^{\#} \in r^{-3} \bar{H}_{\left(-l-\alpha_{2}\right)}^{b}\left(X_{1}^{\#}\right) \tag{6.64}
\end{equation*}
$$

We obtain from (6.62) that

$$
\begin{equation*}
L_{1}^{\# a^{a_{1}}}\left(V_{1}^{\#}, U_{2}^{\#}\right)^{\alpha} u_{2} \in r^{-3} \bar{H}_{\left(-a_{1}-\alpha_{2}\right)}^{b}\left(X_{1}^{\#}\right), \quad a_{1}+\alpha_{1} \leq k \tag{6.65}
\end{equation*}
$$

It follows, in the notation of Appendix B , that (6.64) is equivalent to

$$
\begin{equation*}
L_{1}^{\#^{l}}\left(V_{1}^{\#}, U_{2}^{\#}\right)^{\alpha} u_{2} \in r^{-3} \bar{H}_{\left(-l-\alpha_{2}, \alpha_{2}\right)}^{b}\left(X_{1}^{\#}\right), \quad l+\alpha_{1}+\alpha_{2} \leq k \tag{6.66}
\end{equation*}
$$

Let us observe that $L_{1}^{\#}$ is a second order differential operator for which $\partial X^{\#}=\{Z=$ $2\}$ is non-characteristic. It follows from (6.66) and Proposition B. 1 that by taking $a_{1}=l+\alpha_{2}$

$$
\begin{equation*}
L_{1}^{\# l}\left(V_{1}^{\#}, U_{2}^{\#}\right)^{\alpha} u_{1}^{\#} \in r^{-3} \bar{H}_{(-l)}^{b}\left(X_{1}^{\#}\right), \quad l+|\alpha| \leq k \tag{6.67}
\end{equation*}
$$

Thus the pull-back of (6.67) by $\Psi$ gives (6.61).
The vector fields $r \partial_{r}, U_{2}$ and $\partial_{Z}$ span all smooth vector fields in $\widetilde{X}_{1}$ tangent to $\beta_{1}^{*} \widetilde{R}$, while $r \partial_{r}, U_{2}$ and $3(Z-X) \partial_{X}-Z \partial_{Z}$ span the ones tangent to $\beta_{1}^{*} \widetilde{F}$. Therefore we conclude that the operators in (6.61) span $\Psi_{b}^{1}\left(\widetilde{X}_{1},{ }^{b} N^{*} \beta_{1}^{*} \widetilde{F} \sqcup^{b} N^{*} \beta_{1}^{*} \widetilde{R}\right)$.
Another simple computation shows that

$$
\begin{gather*}
(Z-X) U_{2} \partial_{Z}-\left(1-X^{2}\right) L_{1}=-\left(2 X+Z^{3}-3 Z\right) \partial_{Z}^{2}  \tag{6.68}\\
\left(2-X Z-Z^{2}\right) L_{1}+Z U_{2} \partial_{Z}=-3\left(2 X+Z^{3}-3 Z\right) \partial_{Z} \partial_{X}
\end{gather*}
$$

Since $\partial_{Z}$ is transversal to $\partial X_{1}$, the same method as the one used in the proof of (6.61) shows that if $U_{3}=\left(2 X+Z^{3}-3 Z\right) \partial_{Z}$, then

$$
\begin{equation*}
\left(r \partial_{r}, U_{2}, U_{3}\right)^{\alpha} u_{1} \in r^{-3} L^{2}\left(X_{1}\right), \quad|\alpha| \leq k \tag{6.69}
\end{equation*}
$$

The vector fields $r \partial_{r}, U_{2}$ and $U_{3}$ span all the vector fields tangent to $\beta_{1}^{*} \widetilde{R}$ and $\beta_{1}^{*} \widetilde{F}$. We conclude that away from $\beta_{1}^{*} \widetilde{F}$ and $\beta_{1}^{*} \widetilde{R}, u_{1} \in r^{-3} \bar{H}_{(k)}^{b}\left(X_{1}\right)$ and the extension is trivial. Therefore we may restrict our analysis to a neighbourhood of the intersections of the hypersurfaces and $\partial X_{1}$.
Let us first consider the case where $\beta_{1}^{*} \widetilde{R}$ intersects $\partial X_{1}$ away from $\beta_{1}^{*} \widetilde{F}$. Since this intersection is transversal one can introduce local coordinates $x^{\prime}=X \mp 1, z^{\prime}=Z \mp r \phi_{i}, i=2,3$ in which

$$
\beta_{1}^{*} \widetilde{R}=\left\{x^{\prime}=0\right\}, \quad \partial X_{1}=\left\{z^{\prime}=0\right\}
$$

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Since away from $\beta_{1}^{*} \widetilde{F}$ the vector fields in (6.69) span all the vector fields tangent to $\beta_{1}^{*} \widetilde{R}$, we find that

$$
\begin{equation*}
\left(r \partial_{r}, x^{\prime} \partial_{x^{\prime}}, \partial_{z^{\prime}}\right)^{\alpha} \tilde{u}_{1} \in r^{-3} L^{2}\left(X_{1}\right), \quad|\alpha| \leq k \tag{6.70}
\end{equation*}
$$

Now we deduce from Proposition 6.6 that there exists $\tilde{u}_{1} \in r^{-3} L_{b}^{2}\left(\widetilde{X}_{1}\right)$ such that

$$
\begin{equation*}
\tilde{u}_{1}=u_{1} \text { in } X_{1} \text { and }\left(r \partial_{r}, x^{\prime} \partial_{x^{\prime}}, \partial_{z^{\prime}}\right)^{\alpha} \tilde{u}_{1} \in r^{-3} L_{b}^{2}\left(\tilde{X}_{1}\right), \quad|\alpha| \leq k . \tag{6.71}
\end{equation*}
$$

Finally we analyze the case where $\beta_{1}^{*} \widetilde{F}$ and $\beta_{1}^{*} \widetilde{R}$ intersect the boundary simultaneously. Let us concentrate on the part of the boundary given by $Z=2+r \phi_{3}$ as the other case is analogous. Consider the change of variables

$$
x^{\prime}=2 \frac{X+1}{(Z+1)^{2}}, \quad z^{\prime}=Z-2
$$

which is smooth since $Z+1>0$ near the intersection in question. In these coordinates

$$
\begin{gather*}
\beta_{1}^{*} \widetilde{F}=\left\{x^{\prime}=-z^{\prime}\right\}, \quad \beta_{1}^{*} \widetilde{R}=\left\{x^{\prime}=0\right\}, \text { and } \\
\partial X_{1}=\left\{z^{\prime}=r \phi_{4}\right\}, \quad \phi_{4} \in C^{\infty}\left(\widetilde{X}_{1}\right) . \tag{6.72}
\end{gather*}
$$

From the fact that $r \partial_{r}, U_{2}$ and $U_{3}$ span all the vector fields tangent to $\beta_{1}^{*} \widetilde{R}$ and $\beta_{1}^{*} \widetilde{F}$ we deduce that

$$
\begin{equation*}
\left(r \partial_{r},\left(x^{\prime}+z^{\prime}\right) \partial_{z^{\prime}}, x^{\prime}\left(\partial_{x^{\prime}}-\partial_{z^{\prime}}\right)\right)^{\alpha} u_{1} \in r^{-3} L_{b}^{2}\left(\widetilde{X}_{1}\right), \quad|\alpha| \leq k \tag{6.73}
\end{equation*}
$$

Then we blow-up the submanifold $\left\{x^{\prime}=z^{\prime}=r^{\prime}=0\right\}$ with homogeneity $1-1-1$. Homogeneity of the vector fields in (6.73) and the fact that $\partial X_{1}$ is non-characteristic for $\partial_{z^{\prime}}$ give that the vector fields in (6.73) lift under the $1-1-1$ blow-up to smooth vector fields tangent to the lifts of the hypersurfaces $\beta_{1}^{*} \widetilde{R}$ and $\beta_{1}^{*} \widetilde{F}$ and transversal to the lift of $\partial X_{1}$. Now we can use Proposition 6.6 to extend the lift of $u_{1}$ across the lift of $\partial X_{1}$ to be conormal to the lifts of $\beta_{1}^{*} \widetilde{R}$ and $\beta_{1}^{*} \widetilde{F}$. Then we can show that this in fact gives an extension of $u_{1}$ across $\partial X_{1}$ into the conormal space to the hypersurfaces $\beta_{1}^{*} \widetilde{R}$ and $\beta_{1}^{*} \widetilde{F}$.
We still need to construct the extension map in the region where projective coordinates $r=\left|x_{3}\right|^{\frac{1}{3}} Y=x_{2} / r^{2}, Z=x_{1} / r$ are used. In these coordinates

$$
\begin{gather*}
\beta_{1}^{*} \widetilde{R}=\left\{Y^{3}-1=0\right\}, \quad \beta_{1}^{*} \widetilde{F}=\left\{ \pm 2+Z^{3}-3 Y Z=0\right\}, \\
\partial X_{1}=\left\{Y=\frac{1}{4} Z^{2}+r \phi_{1}\right\}, \quad \phi_{1} \in C^{\infty}\left(\widetilde{X}_{1}\right) . \tag{6.74}
\end{gather*}
$$

Observe that $Y>0$ near $\beta_{1}^{*} \widetilde{R}$ and near the intersection of $\beta_{1}^{*} \widetilde{F}$ and $\partial X_{1}$. Therefore coordinates $(r, X, Z)$ can be used there and the extension can be constructed as above. Away from the two hypersurfaces $u_{1} \in r^{-3} \bar{H}_{(k)}^{b}\left(X_{1}\right)$ and the extension is trivial.
Similarly in the region where coordinates $r=\left|x_{1}\right|, Y=x_{2} / r^{2}, X=x_{3} / r^{3}$ are valid, the intersections of the boundary and the hypersurfaces are contained in the
region where $|Y|>0$ and the extension is constructed as above. Away from the hypersurfaces $u_{1} \in r^{-3} \bar{H}_{(k)}^{b}\left(X_{1}\right)$ and the extension is trivial.
Let us remark again that the linearity and continuity of $E_{2}$ follow from those of the map $S$ of Proposition 6.6 This concludes the proof of Proposition 6.8.

## 7. ESTIMATES FOR THE DIRICHLET PROBLEM

7.1. To solve the mixed problem with the Dirichlet boundary condition we proceed by solving

$$
P u=0 \quad \text { in } \quad X, \quad u \upharpoonright_{\partial X}=f, \quad u \upharpoonright_{X_{-}}=0
$$

where $f$ comes from restricting the solution of the free equation obtained in Theorems 3.9 and 5.6 Thus we first need to characterize $f$.

Let $(x, y)$ be the coordinates in Proposition 3.3 (i) which were used in the definition of $\widetilde{X}_{4}$. Hence

$$
\partial X=\{\rho=0\} \cap \tilde{X}, \quad \Gamma=D \cap \partial X, \quad d \rho \upharpoonright_{\Gamma}=d x_{2}
$$

and $\left(x_{1}, x_{3}, y\right)$ gives a coordinate system on $\partial X$. In $X_{r}$ defined in the beginning of chapter 6 , the 1-2-3 blow-up followed by the $2-1-1$ and $6-3-1$ blow-ups induce a $1-3$ blow-up followed by $2-1$ and $6-1$ blow-ups on the boundary $\beta_{r}^{*} \partial X$. Thus, let us write $Y=\partial X$ and define
$Y_{1}=(Y \backslash \Gamma) \sqcup\left(\mathbb{S}_{1-3}^{1} \times \mathbb{R}^{n-3}\right) \simeq \mathbb{R}_{+} \times \mathbb{S}_{1-3}^{1} \times \mathbb{R}^{n-3}, \quad \mathbb{S}_{1-3}^{1}=\left\{\omega \in \mathbb{R}^{2}: \omega_{1}^{12}+\omega_{2}^{4}=1\right\}$
with the $C^{\infty}$-structure given by the second identification. The blow-down map $\beta_{1}^{\partial}: Y_{1} \rightarrow Y$ is given by

$$
\beta_{1}^{\partial}:(r, \omega, y) \longmapsto\left(r \omega_{1}, r^{3} \omega_{2}, y\right) .
$$

We then define $Y_{2}$ similarly to $\widetilde{X}_{2}$ by blowing-up $\beta_{1}^{\partial^{*}}\left\{x_{1}=0, x_{3} \geq 0\right\} \cap \partial Y_{1}=\Gamma_{1}$ :

$$
Y_{2}=\left(Y_{1} \backslash \Gamma_{1}\right) \sqcup\left(\mathbb{S}_{2-1_{+}}^{1} \times \mathbb{R}^{n-3}\right), \quad \mathbb{S}_{2-1_{+}}^{1}=\left\{\omega \in \mathbb{R}^{2}: \omega_{1}^{2}+\omega_{2}^{4}=1, \omega_{2} \geq 0\right\}
$$

with

$$
Y_{2} \xrightarrow{\beta_{12}^{\partial}} Y_{1} \xrightarrow{\beta_{1}^{\partial}} Y, \quad \beta^{\partial}=\beta_{1}^{\partial} \circ \beta_{12}^{\partial}
$$

where $\beta_{12}^{\partial}(\rho, \omega, y)=\left(\rho^{2} \omega_{1}, \rho \omega_{2}, y\right)$ and where the coordinates in $Y_{1}$ near $\Gamma_{1}$ are chosen so that $\beta_{1}^{\partial}\left(X_{1}, r, y\right)=\left(r X_{1}, r^{3}, y\right)$.

Finally, we have $Y_{3} \xrightarrow{\beta_{23}^{\partial}} Y_{2}$,

$$
Y_{3}=\left(Y_{2} \backslash\left(\beta_{2}^{\partial^{*}}\left\{x_{1}=0, x_{3} \geq 0\right\} \cap \partial Y_{2}\right)\right) \sqcup\left(\mathbb{S}_{6-1_{+}} \times \mathbb{R}^{n-3}\right)
$$

defined as a $C^{\infty}$ manifold with corners in a similar way.
We recall $[28,30]$ that a diffeomorphism between manifolds with corners is a homeomorphism which, together with its inverse, induces $C^{\infty}$ maps on all the boundary faces. The induced blow-up of the boundary is made precise in the following

Lemma 7.1. - There exists a diffeomorphism $f: \beta_{r}^{*} \partial X \rightarrow Y_{3}$ such that $\beta_{r} \circ \imath_{1}=$ $\beta^{\partial} \circ f$, where $\imath_{1}: \beta_{r}^{*} \partial X \hookrightarrow \widetilde{X}_{r}$.

Proof. - The boundary blow-up was defined using the coordinates $\left(x_{1}, x_{3}, y\right)$ used also in the definition of $\widetilde{X}_{r}$. Thus it suffices to check that, after normalizations according to homogeneity, $\beta_{r}^{*} x_{1}, \beta_{r}^{*} x_{3}, \beta_{r}^{*} y$ give coordinates on $\beta_{r}^{*} \partial X$. Since $d \rho \upharpoonright_{\Gamma}=$ $d x_{2}$ this is immediately verified in each projective coordinate system for $\widetilde{X}_{r}$.

Hence we can identify $Y_{3}$ and $\beta_{r}^{*} \partial X$ so that the diffeomorphism $f$ will be omitted below.

This suggests the definition

$$
\begin{equation*}
\mathfrak{S}_{\partial}=\beta^{\partial *}\left(\widetilde{S}_{+} \cap \partial X\right) \sqcup \beta^{\partial *}(\widetilde{F} \cap \partial X) \sqcup \beta^{\partial *}((\widetilde{R} \backslash \widetilde{F}) \cap \partial X) \sqcup \beta^{\partial *}(H \cap \partial X) \sqcup \partial \widetilde{X}_{4} \tag{7.1}
\end{equation*}
$$

where we note that $\mathcal{S}_{\partial}=\imath_{1}^{*} \beta_{2 r}^{*}\left(\beta_{24}\right)_{*} \mathfrak{\delta}$, with $\mathfrak{S}$ defined by (3.7) and appearing in Definition 3.5 of $J_{k} L^{2}(\tilde{X}, H)$.

Recalling that the conormal spaces with non-integral orders of regularity are defined by complex interpolation we can state the restriction result:

Proposition 7.2. - If $u \in J_{k}^{1} L_{c}(\widetilde{X}, H), 1 \leq j \leq n$, then

$$
u \varsigma_{\partial X} \in \beta_{*}^{\partial} I_{k+1 / 2} L_{\nu_{\partial}}^{2}\left(Y_{3}, \mathscr{C}\left(Y_{3}, \mathscr{S}_{\partial}^{+}\right)\right)
$$

where

$$
\begin{align*}
\mathfrak{s}_{\partial}^{+}= & \partial Y_{3} \sqcup \beta^{\partial *}\left(\widetilde{S}_{+} \cap \partial X\right) \sqcup \beta^{\partial *}(\widetilde{F} \cap \partial X) \sqcup  \tag{7.2}\\
& \beta^{\partial *}(((\widetilde{R} \backslash \widetilde{F}) \backslash K) \cap \partial X) \sqcup \beta^{\partial *}((H \backslash K) \cap \partial X),
\end{align*}
$$

$\beta_{*}^{\partial} \nu_{\partial}=d x_{1} d x_{3} d y$, and $K$ in Definition 5.1, with $\epsilon$ chosen as in
The geometry of $\mathfrak{S}_{\partial}^{+}$is shown in Fig. 7.1. The proof will follow easily from the following

Lemma 7.3. - If $\beta^{*} u, \beta^{*} D_{j} u \in I_{k} L_{\nu}^{2}\left(\widetilde{X}_{4}, \widetilde{( }\left(\widetilde{X}_{4}, \mathcal{S}\right)\right)$ then, with $S_{\partial}$ is given by (7.1),

$$
\beta_{r}^{*} u \upharpoonright_{\beta_{r}^{*} \partial X} \in I_{k+1 / 2} L_{\nu_{\partial}}^{2}\left(Y_{3}, \circlearrowright\left(Y_{3}, \varsigma_{\partial}\right)\right) .
$$

Proof. - Let us recall the definition (3.9)

$$
\mathcal{S}=\beta^{*} \widetilde{F} \sqcup \beta^{*} \widetilde{R} \sqcup \beta^{*} \widetilde{R} \sqcup \beta^{*} \widetilde{S}_{+} \sqcup \beta^{*} H \sqcup \beta^{*}(\widetilde{F} \cap \widetilde{R} \backslash B) \sqcup \beta^{*}\left(\widetilde{S}_{+} \cap \widetilde{R} \backslash B\right)
$$



Figure 7.1. The blow-ups on the boundary and $\mathfrak{S}_{\partial}^{+}$in $Y_{3}$
and introduce the notation

$$
\mathcal{S}_{r}=\beta_{2 r}^{*} \circ\left(\beta_{24}\right)_{*} \mathfrak{S}, \quad \mathfrak{S}_{\partial}=\imath_{1}^{*} \mathfrak{S}_{r}
$$

Since the only surfaces intersecting $\beta^{*} D$ are $\beta^{*} \widetilde{S}$ and $\beta^{*} H$, Proposition 6.4 shows that

$$
\left.\beta^{*} v \in I_{k} L_{\nu}^{2}\left(\tilde{X}_{4}, \vartheta\right)\left(\tilde{X}_{4}, \mathfrak{\Im}\right)\right) \Longrightarrow \beta_{a}^{*} v \in I_{k} L_{\nu_{a}}^{2}\left(\tilde{X}_{a}, \vartheta\left(\tilde{X}_{a}, \beta_{r a}^{*} \varsigma_{r}\right)\right)
$$

Since in a neighbourhood of $\beta_{r}^{*} \partial X=\beta_{a}^{*} \partial X, \widetilde{X}_{a}$ and $\widetilde{X}_{r}$ are equal, the assumptions of the lemma imply that

$$
\chi \beta_{r}^{*} u \in I_{k} H_{(1), \nu_{r}}^{b}\left(\widetilde{X}_{r}, \mathscr{S}_{r}\right)
$$

for some $\chi \in C^{\infty}\left(\widetilde{X}_{r}\right), \chi \equiv 1$ near $\beta_{r}^{*} \partial X$.
The surfaces $\beta_{r}^{*} \widetilde{S}_{+}, \beta_{r}^{*} H$ and $\beta_{r}^{*} \widetilde{F}, \beta_{r}^{*} \widetilde{R}$ away from $\beta_{r}^{*}(\widetilde{F} \cap \widetilde{R} \backslash B)$, intersect $\beta_{r}^{*} \partial X$ transversally and are separated from each other. Thus, away from $\beta_{r}^{*}(\widetilde{F} \cap \widetilde{R} \backslash B)$ (see Fig. 3.1), $V \in \mathcal{O}\left(Y_{3}, \mathscr{\varsigma}_{\partial}\right)$ extends to a vector field $\widetilde{V} \in \mathcal{O}\left(\mathscr{(}_{r}\right), \widetilde{V} \upharpoonright_{\beta_{r}^{*} \partial X}=V$, so that

$$
V^{l}\left(u \upharpoonright_{\beta_{r}^{*} \partial X}\right)=\left(\tilde{V}^{l} u\right) \upharpoonright_{\beta_{r}^{*} \partial X} \in H_{\left(\frac{1}{2}\right), \nu_{\partial}}^{b}\left(Y_{3}\right), l \leq k
$$

as $Y_{3} \simeq \beta_{r}^{*} \partial X \subset \widetilde{X}_{r}$ and $\widetilde{V}^{l} u \in H_{(1), \nu_{r}}^{b}\left(\widetilde{X}_{r}\right)$. If $U$ is a sufficiently small neighbourhood of $\beta_{r}^{*}(\widetilde{F} \cap \widetilde{R} \backslash B)$ in $\widetilde{X}_{r}$ then

$$
\mathcal{S}_{r} \upharpoonright_{U}=\left(\beta_{r}^{*} \widetilde{R} \sqcup \beta_{r}^{*} \widetilde{F} \sqcup \beta_{r}^{*}(\widetilde{F} \cap \widetilde{R} \backslash B)\right) \upharpoonright_{U} .
$$

By Proposition 3.3 (see Fig. 3.1) $\beta_{r}^{*} \widetilde{F}, \beta_{r}^{*} \widetilde{R}$ and $\beta_{r}^{*} \partial X$ intersect transversally at $\beta_{r}^{*} \tilde{F} \cap \beta_{r}^{*} \partial X$. We conclude that again $V \in{ }^{*}\left(U \cap \beta_{r}^{*} \partial X, \delta_{\partial} \upharpoonright_{U}\right)$ extends to $\tilde{V} \in \mathcal{O}\left(\widetilde{X}_{r}, \mathscr{S}_{r}\right)$, so that the previous argument is applicable. Thus, it follows by induction that

$$
u \upharpoonright_{\beta_{r}^{*} \partial X} \in I_{k} H_{\left(\frac{1}{2}\right), \nu_{\partial}}^{b}\left(Y_{3}, \stackrel{\varsigma}{\partial} \partial\right) .
$$

We observe that

and

$$
u \in I_{k} H_{(s), \nu_{\partial}}^{b}\left(Y_{3}, \mathcal{O}\left(\mathfrak{s}^{2}\right)\right) \Longleftrightarrow(I+A)^{s} u \in I_{k} L_{\nu_{\partial}}^{2}\left(Y_{3}, \mathscr{o}\left(\mathfrak{s}_{\partial}\right)\right),
$$

if $A \in \Psi_{b}^{1}\left(Y_{3}\right)$ is elliptic and $A \geq 0$. Hence we have

$$
\begin{aligned}
& I_{k} H_{\left(\frac{1}{2}\right)}^{b}\left(Y_{3}, \mathscr{O}\left(\mathcal{S}_{\partial}\right)\right)=\left[I_{k} H_{(1), \nu_{\partial}}^{b}\left(Y_{3}, \mathscr{O}\left(\mathcal{S}_{\partial}\right)\right), I_{k} L_{\nu_{\rho}}^{2}\left(Y_{3}, \mathcal{O}\left(\mathcal{S}_{\partial}\right)\right)\right]_{\frac{1}{2}} \subset \\
& {\left[I_{k+1} L_{\nu_{\partial}}^{2}\left(Y_{3}, \mathcal{O}\left(\mathfrak{\delta}_{\partial}\right)\right), I_{k} L_{\nu_{\partial}}^{2}\left(Y_{3}, \mathcal{O}\left(\varsigma_{\partial}\right)\right)\right]_{\frac{1}{2}}=} \\
& I_{k+\frac{1}{2}} L_{\nu_{\partial}}^{2}\left(Y_{3}, \sigma_{(O)}^{\left.\left(\delta_{\partial}\right)\right)}\right. \text {. }
\end{aligned}
$$

To complete the proof of Proposition 7.2 we observe that (see also Fig. 7.1)

$$
\mathfrak{s}_{\partial}^{+} \upharpoonright_{Y_{3} \backslash \beta^{\partial *}(K \cap \partial X)}=\mathcal{s}_{\partial} \upharpoonright_{Y_{3} \backslash \beta^{\partial *}(K \cap \partial X)}, \quad \delta_{\partial}^{+} \upharpoonright_{\beta^{\partial *}(K \cap \partial X)}=\beta^{\partial *}\left(\widetilde{F}_{-} \cap \partial X\right) \upharpoonright_{\beta^{\partial *}(K \cap \partial X)},
$$

Proceeding as in the proof of Lemma 7.3 shows that if $u_{1} \upharpoonright_{K_{1}} \in I_{k} H_{(1), \nu_{r}}^{b}\left(\tilde{X}_{r},\left(\beta_{r}^{*} \tilde{F} \sqcup\right.\right.$ $\left.\left.\partial \tilde{X}_{r} \upharpoonright_{K_{1}}\right)\right), K_{1}=\beta_{r}^{*} K$, then

$$
\left(u \upharpoonright_{\beta_{r}^{*} \partial X}\right) \upharpoonright_{\beta^{\partial *}(K \cap \partial X)} \in I_{k+\frac{1}{2}} L_{\nu_{\partial}}^{2}\left(\beta^{\partial *}(K \cap \partial X), \beta^{\partial *}(\widetilde{F} \cap \partial X) \sqcup \partial Y_{3}\right) .
$$

Combined with Lemma 7.3 these observations give $u \bigcap_{\partial X} \in \beta_{*}^{\partial} I_{k+\frac{1}{2}} L_{\nu_{\partial}}^{2}\left(Y_{3}, \mathcal{O}_{( }\left(Y_{3}, \mathcal{E}_{\partial}^{+}\right)\right)$.
Proposition 7.2 motivates the following definition:
Definition 7.4. - The conormal space on the boundary, $J_{s} L_{c}^{2}(\partial X, H)$, is defined as

$$
J_{s} L_{c}^{2}(\partial X, H)=\beta_{*}^{\partial} I_{s} L_{\nu_{\partial}}^{2}\left(Y_{3}, \mathcal{O}\left(Y_{3}, \Im_{\partial}^{+}\right)\right)
$$

where the variety $\mathfrak{s}_{\partial}^{+}$is given by (7.2).

A slightly stronger formulation of the restriction result is now given in
Theorem 7.5. - If $u \in J_{s}^{1} L^{2}(\tilde{X}, H)$ then $u \upharpoonright_{\partial X} \in J_{s+\frac{1}{2}} L^{2}(\partial X, H)$.
Proof. - The conormal space in the definition of $J_{t} L^{2}(\partial X, H)$ is an interpolations space - see Appendix B. Since for non-integral orders $J_{s}^{1} L^{2}(\widetilde{X}, H)$ is defined by interpolation between $[s]$ and $[s]+1$ the theorem follows from Proposition 7.2.

The main result of this section is
Theorem 7.6. - If $u \in L_{\mathrm{loc}}^{2}(X)$ is the solution of

$$
\begin{equation*}
P u=0 \quad \text { in } \quad X, \quad u \upharpoonright_{\partial X}=f \in J_{s} L_{c}^{2}(\partial X, H), \quad u \upharpoonright_{X_{-}}=0, \quad f \upharpoonright_{(\partial X)_{-}}=0 \tag{7.4}
\end{equation*}
$$

then there exists $\tilde{u} \in J_{s} L_{\text {loc }}^{2}(\tilde{X}, H)$ such that $u=\tilde{u} \upharpoonright_{X}$.
Theorem 7.1 is proved by finding a microlocally characterized space containing $J_{k} L_{c}^{2}(\partial X, H)$ (Proposition 7.7) and by using microlocal models for the components of that space. The propagation estimates obtained in the Friedlander model (Proposition 7.9) give $u$ in terms of marked Lagrangian spaces on a manifold with boundary (Proposition 7.19). The extension property for those spaces obtained in chapter 6 completes the proof.

Let us start by recalling from Section 2 the defintion of the following smooth Lagrangian submanifolds of $T^{*} \partial X \backslash 0$ :
$\Lambda_{00}=N^{*} \Gamma, \quad \Lambda_{11}=N^{*}(\widetilde{F} \cap \partial X), \quad \Lambda_{31}=N^{*}((\widetilde{R} \backslash \widetilde{F}) \cap \partial X), \quad \Lambda_{13}=N^{*}(H \cap \partial X)$,
and

$$
\Lambda_{21} \cup \Lambda_{23}=N^{*}(\widetilde{S} \cap \partial X), \quad \Lambda_{11} \cap \Lambda_{23}=\emptyset, \quad \Lambda_{13} \cap \Lambda_{21}=\emptyset .
$$

From these we obtain Lagrangians with boundaries:

$$
\Lambda_{i j}^{ \pm}=\Lambda_{i j} \cap\left\{ \pm x_{j} \geq 0\right\}, \quad j=1,3, i=1,2,3 .
$$

in terms of coordinates $(x, y)$ of Proposition 3.1.
The sub- and super-marked Lagrangian spaces introduced in chapter 4 now enter in

Proposition 7.7. - If $J_{k} L_{c}^{2}(\partial X, H)$ is given by Definition 7.4, then

$$
\begin{array}{r}
J_{k} L_{c}^{2}(\partial X, H) \subset \sum_{i=1,3} I_{k} L^{2}\left(\partial X ; \Lambda_{00}, J_{\Lambda_{00} \cap \Lambda_{2 i}}^{1} \Lambda_{2 i}, 2\right)+\sum_{i=1,3} I_{k} L^{2}\left(\partial X ; \Lambda_{2 i}^{+}, 2\right)+ \\
I_{k} L^{2}\left(\partial X ; \Lambda_{3 i}^{+}, 2\right)+I_{k} L^{2}\left(\partial X ; \Lambda_{13}^{+}, 2\right)+\sum_{i=1}^{2} I_{k} L^{2}\left(\partial X ; \Lambda_{11}^{\operatorname{sgn}\left((-1)^{i}\right)}, 2\right) \tag{7.6}
\end{array}
$$

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Proof. - We need to consider another blow-up $\beta_{34}^{\partial}: Y_{4} \longrightarrow Y_{3}, \beta_{4}^{\partial}=\beta^{\partial} \circ \beta_{34}^{\partial}$, obtained by successive 2-1 and 6-1 blow-ups of $\beta^{\partial *}\left(\left\{x_{1}=0, x_{3} \leq 0\right\}\right) \cap \partial Y_{3}$ identical to the blow-ups used in the construction of $Y_{3}$.
We easily see that

$$
\beta_{4 *}^{\partial} I_{k} L_{\nu_{\partial 4}}^{2}\left(Y_{4} ; ণ\left(\beta_{34}^{\partial *} \mathfrak{S}_{\partial}^{+} \sqcup \partial Y_{4}\right)\right) \supset J_{k} L_{c}^{2}(\partial X, H), \quad \beta_{4 *}^{\partial} \nu_{\partial 4}=d x_{1} d x_{3} d y
$$

and we shall prove that the left hand side is contained in the right hand side of (7.6). We first prove that


Figure 7.2. The geometry in $Y_{4}$

$$
\begin{gather*}
I_{k} L_{\nu_{\partial 4}}^{2}\left(Y_{4} ; \gamma\left(\beta_{4}^{\partial *}(\tilde{F} \cap \partial X) \sqcup \partial Y_{4}\right)\right) \xrightarrow{\beta_{4 *}} \\
\sum_{i=1,3} I_{k} L^{2}\left(\partial X ; \Lambda_{00}, J_{\Lambda_{00} \cap \Lambda_{2 i}}^{1} \Lambda_{2 i}, 2\right)+I_{k} L^{2}\left(\partial X ; \Lambda_{11}, 2\right) \tag{7.7}
\end{gather*}
$$

In particular this gives

$$
\begin{equation*}
I_{k} L_{\nu_{\partial 4}}^{2}\left(Y_{4} ; \partial Y_{4}\right) \xrightarrow{\beta_{4 *}^{*}} \sum_{i=1,3} I_{k} L^{2}\left(\partial X ; \Lambda_{00}, J_{\Lambda_{00} \cap \Lambda_{2 i}}^{1} \Lambda_{2 i}, 2\right)+I_{k} L^{2}\left(\partial X ; \Lambda_{11}, 2\right) . \tag{7.8}
\end{equation*}
$$

We note that in view of Proposition 4.3 we can replace the second term on the right hand side by the last term in (7.6). To prove (7.7) we first observe that we can
change coordinates in $\partial X$ such that $\widetilde{F} \cap \partial X=\left\{x_{3}=0\right\}, \Gamma=\left\{x_{1}=x_{3}=0\right\}$ and $J_{\Lambda_{00} \cap \Lambda_{23}}^{1} \Lambda_{23}=J_{\Lambda_{00} \cap \Lambda_{23}}^{1}\left(N^{*}\left\{x_{1}=0\right\}\right)$ and that this change of coordinates lifts to a smooth diffeomorphism on $Y_{4}$. To construct such a change of variables we let ( $x_{1}, x_{3}$ ) be smooth coordinates near $\Gamma$ in which $H \cap \partial X=\left\{x_{1}=0\right\}$. Since $\widetilde{F} \cap \partial X$ is second order tangent to $H \cap \partial X$ at $\Gamma$, it follows that $\widetilde{F} \cap \partial X=\left\{x_{3}+f(x, y) x_{1}^{3}=0\right\}$, with $f(x, y) \neq 0$. For simplicity we assume that $f(0,0)=1$. Now a direct calculation shows that there exists a smooth function $A\left(s, x_{1}, x_{3}\right)$ for $s \in[0,1]$ and $x_{1}, x_{3}$ small, such that

$$
A\left(s, x_{1}, x_{3}\right) x_{1} \partial_{x_{1}}\left(x_{3}+x_{1}^{3}+s(f(x, y)-1) x_{1}^{3}\right)=-(f(x, y)-1) x_{1}^{3}
$$

Thus, as in section 3, we obtain, by integrating the vector field $V_{s}=A\left(s, x_{1}, x_{3}\right) x_{3} \partial_{x_{3}}$, a one parameter family smooth diffeomorphisms $\phi_{s}$ fixing $\left\{x_{1}=0\right\}$ and satisfying

$$
\phi_{s}^{*}\left(x_{3}+x_{1}^{3}+s(f(x, y)-1) x_{1}^{3}\right)=x_{3}+x_{1}^{3}
$$

In particular the map $\phi_{1}$ fixes $\left\{x_{1}=0\right\}$ and $\phi_{1}^{*}(\widetilde{F} \cap \partial X)=\left\{x_{3}+x_{1}^{3}=0\right\}$. We are going to show below that the vector field $x_{1} \partial_{x_{1}}$ lifts to a smooth vector field in $Y_{4}$ which is tangent to $\partial Y_{4}$, therefore it follows that $\phi_{1}$ coordinates lifts to a smooth diffeomorphism on $Y_{4}$ preserving $\partial Y_{4}$.

Proposition 4.3 shows that the right hand side of (7.7) is equal to

$$
I_{k} L^{2}\left(\partial X, \Lambda_{00} \sqcup \Lambda_{11}\right)+I_{k} L^{2}\left(\partial X, \Lambda_{00}, J_{\Lambda_{00} \cap \Lambda_{23}}^{1} \Lambda_{23}, 2\right)
$$

Thus Definition 4.5, a microlocal partition of unity and the proof of Lemma 4.2 show that

$$
\begin{array}{r}
u \in I_{k} L^{2}\left(\partial X, \Lambda_{00} \sqcup \Lambda_{11}\right)+I_{k} L^{2}\left(\partial X, \Lambda_{00}, J_{\Lambda_{00} \cap \Lambda_{23}}^{1} \Lambda_{23}, 2\right) \Longleftrightarrow  \tag{7.9}\\
\left(x_{3} D_{x_{3}}\right)^{k_{1}}\left(x_{1} D_{x_{1}}\right)^{k_{2}}\left(x_{3}^{3} D_{x_{1}}\right)^{k_{3}} u \in L^{2}(\partial X), k_{1}+k_{2}+k_{3} \leq k .
\end{array}
$$

To prove that (7.7) holds one needs to show that for $u \in I_{k} L_{\nu_{\partial 4}}^{2}\left(Y_{4} ; \widetilde{\sigma}\right)\left(\beta_{4}^{\partial *}(\widetilde{F} \cap \partial X) \sqcup\right.$ $\left.\partial Y_{4}\right)$ ) and $v=\beta_{4 *} u$ one has

$$
\begin{equation*}
\left(x_{3} D_{x_{3}}\right)^{k_{1}}\left(x_{1} D_{x_{1}}\right)^{k_{2}}\left(x_{3}^{3} D_{x_{1}}\right)^{k_{3}} v \in L^{2}(\partial X), k_{1}+k_{2}+k_{3} \leq k \tag{7.10}
\end{equation*}
$$

Since the operators in (7.10) are smooth vector fields, (7.7) is a consequence of

$$
\begin{equation*}
\beta_{4}^{\partial^{*}}\left(x_{3} \partial_{x_{3}}, x_{1} \partial_{x_{1}}, x_{3}^{3} \partial_{x_{1}}\right) \subset \vartheta\left(\beta_{4}^{\partial^{*}}(\widetilde{F} \cap \partial X) \sqcup \partial Y_{4}\right) . \tag{7.11}
\end{equation*}
$$

In other words one needs to show that the vector fields in (7.9) lift under $\beta_{4}^{\partial}$ to smooth vector filelds in $Y_{4}$ tangent to $\beta_{4}^{\partial^{*}}(\widetilde{F} \cap \partial X)$. To prove (7.11) we compute the lifts of these vector fields in twelve projective coordinate systems. Consider the first set of projective coordinates

$$
\begin{aligned}
& \left(\rho, X_{1}\right) \mapsto\left(\rho X_{1}, \pm \rho^{3}\right)=\left(x_{1}, x_{3}\right), \\
& \left(\rho, X_{3}\right) \mapsto\left( \pm \rho, \rho^{3} X_{3}\right)=\left(x_{1}, x_{3}\right) .
\end{aligned}
$$

In these coodinates the lift of the vector fields in (7.11) are spanned, over $C^{\infty}\left(Y_{1}\right)$, by

$$
\begin{gather*}
\rho \partial_{\rho}, \quad X_{1} \partial_{X_{1}}, \rho^{8} \partial_{X_{1}},  \tag{7.12}\\
\rho \partial_{\rho}, \quad X_{3} \partial_{X_{3}} .
\end{gather*}
$$

The vector fields $\rho \partial_{\rho}$ and $X_{3} \partial_{X_{3}}$ are tangent to $\beta_{1}^{\partial^{*}}(\widetilde{F} \cap \partial X) \sqcup \partial Y_{1}$ which can be identified with $\beta_{4}^{\partial^{*}}(\widetilde{F} \cap \partial X) \sqcup \partial Y_{4}$ away from $\left\{X_{1}=\rho=0\right\}$.
Next we blow-up $\left\{X_{1}=\rho=0\right\}$ with homogeneity $2-1$. Consider the second set of projective coordinates

$$
\begin{gathered}
\left(R, x_{1}^{\prime}\right) \mapsto\left(R, R^{2} x_{1}^{\prime}\right)=\left(\rho, X_{1}\right) \\
\left(R, \rho^{\prime}\right) \mapsto\left(\rho^{\prime} R^{2}, \pm R^{2}\right)=\left(\rho, X_{1}\right)
\end{gathered}
$$

The lift of the vector fields in (7.12) are spanned by

$$
\begin{gather*}
R \partial_{R}, ; x_{1}^{\prime} \partial_{x_{1}^{\prime}}, R^{6} \partial_{x_{1}^{\prime}}  \tag{7.13}\\
R \partial_{R}, \rho^{\prime} \partial_{\rho^{\prime}}
\end{gather*}
$$

Away from $\left\{x_{1}^{\prime}=R=0\right\} Y_{2}$ can be identified with $Y_{4}$ and one easily sees that the second set of vector fields clearly span $\left(Y_{4}\right)$. Finally one needs to blow-up the submanifold $\left\{x_{1}^{\prime}=R=0\right\}$ with homogeneity $6-1$. Consider coordinates

$$
\begin{gathered}
(r, T) \mapsto\left( \pm T^{6}, r T\right)=\left(x_{1}^{\prime}, R\right) \\
\left(r_{1}, T_{1}\right) \mapsto\left(T_{1}^{6} r_{1}, T_{1}\right)=\left(x_{1}^{\prime}, R\right)
\end{gathered}
$$

It is easy to see that the lift of the vector fields in (7.13) is in the span of

$$
\begin{gathered}
T \partial_{T}, r \partial_{r} \\
\partial_{r_{1}}, T_{1} \partial_{T_{1}}
\end{gathered}
$$

Therefore they are in $\left(\partial Y_{4}\right)$. This concludes the proof of (7.7).
We now want

$$
\begin{gather*}
I_{k} L_{\nu_{\partial 4}}^{2}\left(Y_{4} ; \%\left(\beta_{4}^{\partial *}\left((\widetilde{R} \backslash \widetilde{F}) \cap \partial X \cap\left\{x_{1} \geq 0\right\}\right) \sqcup \partial Y_{4}\right)\right) \xrightarrow{\beta_{4} *} \\
\sum_{i=1,3} I_{k} L^{2}\left(\partial X ; \Lambda_{00}, J_{\Lambda_{00} \cap \Lambda_{2 i}}^{1} \Lambda_{2 i}, 2\right)+I_{k} L^{2}\left(\partial X, \Lambda_{00} \sqcup \Lambda_{11}\right)+  \tag{7.14}\\
I_{k} L^{2}\left(\partial X ; \Lambda_{3}^{+}, 2\right)+I_{k} L^{2}\left(\partial X, \Lambda_{11}, 2\right),
\end{gather*}
$$

To prove (7.14) we proceed by a partition of unity in $Y_{4}$. Let $\phi \in C_{c}^{\infty}\left(Y_{4}\right)$ be such that $\phi=1$ near $\beta_{4}^{\partial^{*}}((\widetilde{R} \backslash \widetilde{F}) \cap \partial X) \cap\left\{\beta_{4}^{\partial^{*}} x_{1} \geq 0\right\}$ and $\phi$ is supported away from the boundary faces introduced by the $2-1$ and $6-1$ blow-ups of $\beta_{3}^{\partial^{*}}\left(\left\{x_{1}=0, x_{3} \leq 0\right\}\right) \cap$ $\partial Y_{3}$. Let $\left.u \in I_{k} L_{\nu_{\partial 4}}^{2}\left(Y_{4} ;{ }^{\circ}\right)\left(\beta_{4}^{\partial *}\left((\widetilde{R} \backslash \widetilde{F}) \cap \partial X \cap\left\{x_{1} \geq 0\right\}\right) \sqcup \partial Y_{4}\right)\right)$. Then by (7.8) $\beta_{4 *}^{\partial}(1-\phi) u \in \sum_{i=1,3} I_{k} L^{2}\left(\partial X ; \Lambda_{00}, J_{\Lambda_{00} \cap \Lambda_{2 i}}^{1} \Lambda_{2 i}, 2\right)+I_{k} L^{2}\left(\partial X ; \Lambda_{11}, 2\right)$.

From the support properties of $\phi$ we deduce that

$$
\beta_{4 *}^{\partial}(\phi u)=\beta_{3 *}^{\partial}(\phi u) .
$$

Now it follows from Lemma 4.2 that

$$
\beta_{4 *}^{\partial}(\phi u) \in I_{k} L_{\mathrm{loc}}^{2}\left(\partial X, \Lambda_{3}^{+}, 2\right)+\beta_{3}^{*} H_{(k)}^{b}\left(\partial Y_{3}\right) .
$$

Thus (7.14) follows from (7.8).
The push forward of the other components of $\beta_{34}^{*} \varsigma_{\partial}^{+}$are handled in a similar way and the details will be left to the reader.

Remark 7.8. - The space $J_{k} L^{2}(\partial X, H)$ is strictly contained in the microlocal space defined in Proposition 7.7. The extra terms are in $I_{k} L^{2}\left(\partial X ; \Lambda_{23}^{+}, 2\right)$ and $I_{k} L^{2}\left(\partial X ; \Lambda_{13}^{+}, 2\right)$. They are explained by the additional blow-ups in $Y_{4}$ needed for the first term in the right hand side of (7.6) and the singularities on $\partial Y_{4} \backslash \beta_{34}^{\partial *} \partial Y_{3}$ produced by $I_{k} L^{2}\left(\partial X ; \Lambda_{23}^{+}, 2\right)$ and $I_{k} L^{2}\left(\partial X ; \Lambda_{13}^{+}, 2\right)$. Otherwise the push-forward is sharp.

To solve (7.4) we use the forward Melrose-Taylor diffractive parametrix [35] which we shall briefly recall:

$$
\begin{equation*}
T=\widetilde{T} \circ L, \quad L \in I^{0}\left(\mathbb{R}^{n-1} \times \partial X, \mathscr{g}^{-1^{\prime}}\right), \quad \widetilde{T}: \mathscr{G}^{\prime}\left(\mathbb{R}^{n-1}\right) \longrightarrow \mathfrak{\vartheta}(X) \tag{7.15}
\end{equation*}
$$

where $P \widetilde{T} v \in C^{\infty}(X), \widetilde{T} v \upharpoonright_{\partial X}=J v, \widetilde{T} v \upharpoonright_{X_{-}}=0, J \in I^{0}\left(\partial X \times \mathbb{R}^{n-1}, \mathscr{g}^{\prime}\right)$ is an elliptic (in an appropriate cone) Fourier Integral Operator and $L$ its microlocal inverse. The space $\mathfrak{g i}(X)$ used here gives, in an invariant way, the functions 'smooth in the direction normal to the boundary' - see [14, 25, 28].

Let $\Gamma \subset T^{*} \mathbb{R}^{n-1} \backslash 0$ be a small conic neighbourhood of $(0 ;(0,1,0, \cdots, 0))$. Then for $f \in \mathcal{E}^{\prime}\left(\mathbb{R}^{n-1}\right), W F(f) \subset \Gamma$,

$$
\begin{equation*}
\widetilde{T} f(z)=\int\left(g(z, \xi) A_{+}(\zeta)+h(z, \xi) A_{+}^{\prime}(\zeta)\right) A_{+}\left(\zeta_{0}\right)^{-1} e^{i \theta(z, \xi)} \hat{f}(\xi) d \xi \tag{7.16}
\end{equation*}
$$

with the phase functions $\zeta, \zeta_{0}=\zeta \upharpoonright_{\partial X \times \mathbb{R}^{n-1}}$ and $\theta$ homogeneous of degree $2 / 3$ and 1 respectively. The amplitudes $g \in S^{0}\left(X ; \mathbb{R}^{n-1}\right)$ and $h \in S^{-1 / 3}\left(X, \mathbb{R}^{n-1}\right)$ are supported in a conic neighbourhood of $\left(z_{0},(0, \cdots, 0,1)\right), z_{0} \in \partial X$ and satisfy appropriate transport equations. Most importantly $g(0,(0,1,0, \cdots, 0))=1$ and $h \upharpoonright_{\partial X \times \mathbb{R}^{n-1}}=0$.

The construction of $\zeta$ and $\theta$ exploits the equivalence of glancing hypersurfaces which we will now discuss (see [25, 35], for detailed presentation and proofs). Let $\widetilde{\Gamma} \subset T^{*} \mathbb{R}^{n} \backslash 0$ be an open cone, $\widetilde{\Gamma} \upharpoonright_{x_{2}=0} \subset \Gamma$, where we denote the coordinates in $\mathbb{R}^{n}$ by $(x, y), x \in \mathbb{R}^{3}, y \in \mathbb{R}^{n-3}$ with $\mathbb{R}^{n-1}=\partial \mathbb{R}_{+}^{n}=\left\{x_{2}=0\right\} \cap \mathbb{R}^{n}$ (compare Proposition 3.1) and let $m \in T_{\partial X}^{*} \tilde{X} \backslash 0$ be a glancing point for $p=\sigma_{2}(P)$. Then there exists a canonical transformation

$$
\tilde{\chi}: \tilde{\Gamma} \longrightarrow T^{*} \tilde{X} \backslash 0, \quad \tilde{\chi}(0 ;(0,1,0, \cdots, 0))
$$

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such that

$$
\begin{align*}
& \tilde{\chi}\left(\left\{\xi_{2}^{2}-x_{2} \xi_{3}^{2}-\xi_{3} \xi_{1}=0\right\} \cap \widetilde{\Gamma}\right) \subset p^{-1}(0)  \tag{7.17}\\
& \tilde{\chi}\left(\left\{x_{2}=0\right\} \cap \widetilde{\Gamma}\right) \subset T_{\partial X}^{*} \widetilde{X}
\end{align*}
$$

The equivalence $\tilde{\chi}$ induces a canonical transformation on the boundary

$$
\chi_{\partial}: \Gamma \longrightarrow T^{*} \partial X \backslash 0,
$$

and the graph of $\chi_{\partial}, \mathscr{J}$, gives then the canonical relation of the elliptic Fourier Integral Operator $J$. The main geometric property of $\chi_{\partial}$ is the intertwining of the billiard ball maps defined in chapter 2 :

$$
\begin{equation*}
\chi_{\partial}^{-1} \circ \delta^{ \pm} \circ \chi_{\partial}=\delta_{0}^{ \pm} \tag{7.18}
\end{equation*}
$$

where $\delta_{0}^{ \pm}$corresponds to the model glancing hypersurfaces in the left hand side of (7.17).

There is a substantial amount of freedom in choosing $\tilde{\chi}$. Thus, in addition to (7.17) we can also have

$$
\begin{equation*}
\chi_{\partial}\left(\Lambda_{0} \cap \Gamma\right) \subset \Lambda_{00} \quad \text { or } \quad \chi_{\partial}\left(\Lambda_{1} \cap \Gamma\right) \subset \Lambda_{1 i}, \quad i=1 \quad \text { or } 3 \tag{7.19}
\end{equation*}
$$

where $\Lambda_{0}=N^{*}\left\{x_{1}=x_{3}=0\right\} \subset T^{*} \mathbb{R}^{n-1} \backslash 0$ and $\Lambda_{1}=N^{*}\left\{x_{3}+x_{1}^{3} / 3=0\right\} \subset$ $T^{*} \mathbb{R}^{n-1} \backslash 0$. In fact, one can apply Theorem 4.2 .6 of [36] in the same way as in Sect. 3 of [37] and Sect. 3 of [51]. We should note that in (7.19) one needs to choose a different $\chi_{\partial}$ in the case of each Lagrangian $\Lambda_{0}, \Lambda_{1 i}, i=1,3$.

We recall from chapter 2 that

$$
\Lambda_{31}^{ \pm}=\delta^{ \pm} \Lambda_{11}^{ \pm}, \quad \Lambda_{21}^{ \pm} \cup \Lambda_{23}^{ \pm}=\delta^{ \pm} \Lambda_{0}
$$

and

$$
\Lambda_{33}^{ \pm}=\delta^{ \pm} \Lambda_{13}^{ \pm}
$$

where $\Lambda_{33}$ is a Lagrangian simply tangent to $\Lambda_{23}$ at $\Lambda_{23} \cap \Lambda_{13}$. It corresponds to the reflection of the false front $H$.

Thus, thanks to (7.18) in each case of (7.19) we also obtain

$$
\begin{array}{ll}
\chi_{\partial}\left(\Lambda_{2} \cap \Gamma\right) \subset \Lambda_{2 i}, & i=1 \quad \text { or } 3 \\
\chi_{\partial}\left(\Lambda_{3} \cap \Gamma\right) \subset \Lambda_{3 i}, & i=1 \quad \text { or } 3
\end{array}
$$

where $\Lambda_{m}=N^{*}\left\{x_{3}+x_{1}^{3} /\left(3 m^{2}\right)=0\right\}, m \neq 0$.
If the phases in (7.16) are considered formally then the expected wave front set relation for $\widetilde{T}$ becomes $\mathfrak{F}^{+} \circ \mathfrak{C}_{0}^{-}$, where $\mathfrak{C}_{0}$ is the model square root of the billiard ball map (see [36] and (7.26) below) and

$$
\mathfrak{F}^{ \pm}=\left\{\left(z, \phi_{z}^{ \pm} ; \phi_{\xi}^{ \pm}, \xi\right): \phi^{ \pm}(z, \xi)=\theta \mp \frac{2}{3}(-\zeta(z, \xi))^{3 / 2}\right\} \subset T^{*} \tilde{X} \backslash 0 \times T^{*} \mathbb{R}^{n-1} \backslash 0
$$

with $\mathcal{F}=\mathscr{F}^{+} \cup \mathcal{F}^{-}$, a smooth canonical relation generated by $\tilde{\phi}(z, \xi, \tau)=-\xi_{1}^{-2} \tau^{3} / 3-$ $\tau \xi_{1}^{-2 / 3} \zeta(z, \xi)+\theta(z, \xi)$. We also have $\mathcal{F}_{0}=\mathcal{F}_{0}^{+} \cup \mathcal{F}_{0}^{-}$for the model case.

The construction of $\zeta$ and $\theta$ in [35] shows that

$$
\begin{equation*}
\mathcal{F}=\left\{(r, s) \in \tilde{\chi}(\widetilde{\Gamma}) \times \Gamma \subset T^{*} \tilde{X} \backslash 0 \times T^{*} \mathbb{R}^{n-1} \backslash 0:\left(\tilde{\chi}^{-1}(r), s\right) \in \mathcal{F}_{0}\right\} \tag{7.20}
\end{equation*}
$$

and we observe that consequently for each choice of $\chi_{\partial}$ in (7.19)

$$
\begin{aligned}
& \chi_{\partial}\left(\Lambda_{1} \cap \Gamma\right) \subset \Lambda_{11} \Longrightarrow N^{*} \widetilde{F} \cap \widetilde{\chi}(\widetilde{\Gamma}) \subset \mathcal{F}\left(\Lambda_{0} \cap \Gamma\right), N^{*} \widetilde{R} \cap \widetilde{\chi}(\widetilde{\Gamma}) \subset \mathcal{F}\left(\Lambda_{2} \cap \Gamma\right) \\
& \chi_{\partial}\left(\Lambda_{0} \cap \Gamma\right) \subset \Lambda_{00} \Longrightarrow N^{*} \widetilde{S} \cap \widetilde{\chi}(\widetilde{\Gamma}) \subset \mathcal{F}\left(\Lambda_{1} \cap \Gamma\right) \\
& \chi_{\partial}\left(\Lambda_{1} \cap \Gamma\right) \subset \Lambda_{13} \Longrightarrow N^{*} H \cap \widetilde{\chi}(\widetilde{\Gamma}) \subset \mathcal{F}\left(\Lambda_{0} \cap \Gamma\right), N^{*} H_{+} \cap \widetilde{\chi}(\widetilde{\Gamma}) \subset \mathcal{F}\left(\Lambda_{2} \cap \Gamma\right)
\end{aligned}
$$

where $H_{+}$is the reflection ${ }^{1}$ of $H$ by $\partial X$ and where in the second case we consider different choices of $\widetilde{\chi}$ and $\widetilde{\Gamma}$ giving localization near $B$ and $D$ respectively. We note that by choosing different conic neighbourhoods $\Gamma$ and $\widetilde{\Gamma}$ we can reverse the inclusions above.
7.2. In this subsection we shall consider the model problem in

$$
\mathbb{R}_{+}^{3}=\left\{x \in \mathbb{R}^{3}: x_{2}>0\right\}
$$

with $P_{0}=D_{x_{2}}^{2}-x_{2} D_{x_{3}}^{2}-D_{x_{1}} D_{x_{3}}$,

$$
\begin{equation*}
P_{0} u=0 \text { in } \mathbb{R}_{+}^{3}, u \upharpoonright_{\partial \mathbb{R}_{+}^{3}}=u_{0} \in \mathfrak{G}^{\prime}\left(\mathbb{R}^{2}\right), \quad u \upharpoonright_{\phi(x) \ll 0}=0, \tag{7.21}
\end{equation*}
$$

where $\phi(x)$ is the time function for $P_{0}$, near zero chosen to be $x_{1}+x_{3}$. The 1-2-3 homogeneity of the problem has already been stressed in chapter 3 and here we shall use it microlocally near $m_{0}=(0 ;(0,0,1))$.

The solution operator for (7.21) is given explicitly as

$$
\begin{equation*}
u=T_{0} u_{0}, \quad T_{0} u_{0}(x)=\frac{1}{(2 \pi)^{2}} \int \frac{A_{+}(\zeta)}{A_{+}\left(\zeta_{0}\right)} \hat{u}_{0}\left(\xi_{1}, \xi_{3}\right) e^{i\left(\xi_{1} x_{1}+\xi_{3} x_{3}\right)} d \xi_{1} d \xi_{3} \tag{7.22}
\end{equation*}
$$

where

$$
\zeta=-\xi_{3}^{-1 / 3}\left(\xi_{1}+x_{2} \xi_{3}\right), \quad \zeta_{0}=\zeta\left\lceil_{\partial \mathbb{R}_{+}^{3} \times \mathbb{R}^{3}}=-\xi_{3}^{-1 / 3} \xi_{1} .\right.
$$

Let $\Gamma_{0} \subset T^{*} \partial \mathbb{R}_{+}^{3} \backslash 0$ be a small connected open conic neighbourhood of $(0,(0,1))$ and let $\Gamma \subset T^{*} \mathbb{R}^{3} \backslash 0$ be a connected open conic neighbourhood of $(0 ;(0,0,1))$ such that

$$
W F_{b}\left(\chi T u_{0}\right) \subset \jmath_{*} \Gamma \quad \text { if } \quad W F\left(u_{0}\right) \subset \Gamma_{0}, \quad \chi \in C_{0}^{\infty}\left(\mathbb{R}^{3}\right), \quad \text { supp } \chi \text { near } 0
$$

and where $\jmath$ is the natural inclusion $\jmath: T^{*} \mathbb{R}_{+}^{3} \backslash 0 \rightarrow{ }^{b} T^{*} \mathbb{R}_{+}^{3} \backslash 0$. The existence of $\Gamma$ follows from propagation of singularities for the diffractive boundary value problem and in this case can be easily seen directly.

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We also define the microlocal parametrix near $\Gamma_{0}$ :

$$
\begin{align*}
& T_{0}^{\sharp} u_{0}(x)=\frac{1}{(2 \pi)^{2}} \int \beta\left(\xi_{3}\right) \psi\left(\xi_{3}^{-1} \xi_{2}\right) \frac{A_{+}(\zeta)}{A_{+}\left(\zeta_{0}\right)} \widehat{u}_{0}\left(\xi_{1}, \xi_{3}\right) e^{i\left(x_{1} \xi_{1}+x_{3} \xi_{3}\right)} d \xi_{1} d \xi_{3}  \tag{7.23}\\
& \quad \beta \in C^{\infty}([1, \infty)), \beta(t)=1 \quad \text { if } t>2, \quad \psi \in C_{0}^{\infty}(\mathbb{R}), \psi=1 \text { near } 0
\end{align*}
$$

and we are interested in its mapping properties.
To exploit the homogeneity in a systematic way, let us now introduce

$$
\begin{gathered}
Z=[1, \infty) \times \mathbb{R}, \quad Z_{+}=[1, \infty) \times \mathbb{R} \times \mathbb{R}_{+} \\
(\lambda, x) \in Z, \quad(\lambda, x, y) \in Z_{+}
\end{gathered}
$$

and the isometries

$$
W: L^{2}\left(\mathbb{R}^{2}\right) \longrightarrow \lambda^{\frac{1}{6}} L^{2}(Z), \quad W_{+}: L^{2}\left(\mathbb{R}_{+}^{3}\right) \longrightarrow \lambda^{\frac{1}{2}} L^{2}\left(Z_{+}\right)
$$

extending the maps defined for Schwartz function as

$$
\begin{aligned}
u_{0} & \longmapsto W u_{0}(\lambda, x)=\sqrt{2 \pi} \int u_{0}\left(\lambda^{-\frac{1}{3}} x, x_{3}\right) e^{-i \lambda x_{3}} d x_{3}, \\
u & \longmapsto
\end{aligned} W_{+} u(\lambda, x, y)=2 \pi \int u\left(\lambda^{-\frac{1}{3}} x, \lambda^{-\frac{2}{3}} y, x_{3}\right) e^{-i \lambda x_{3}} d x_{3} .
$$

This corresponds to a non-homogeneous blow-up on the Fourier transfer side with

$$
\left(\lambda^{\frac{1}{3}}, \xi, \eta\right), \quad \xi \text { dual to } x, \eta \text { dual to } y
$$

giving the projective coordinate near the lift of $(0,0,1)$. We observe that because of the cut-off $\beta$ in (7.23), there exists a unique operator

$$
S: \delta(Z) \longrightarrow C^{\infty}\left(Z_{+}\right)
$$

such that $S W=W_{+} T^{\sharp}$. The formula for $T^{\sharp}$ provides an explicit expression for $S$ :

$$
\begin{equation*}
S v(\lambda, x, y)=\frac{1}{(2 \pi)^{\frac{3}{2}}} \int \psi\left(\lambda^{-\frac{2}{3}} \xi\right) \beta(\lambda) \frac{A_{+}(-\xi-y)}{A_{+}(-\xi)} \widehat{v}(\xi) e^{i x \xi} d \xi \tag{7.24}
\end{equation*}
$$

and since our considerations are local in the original coordinates $\left(x_{1}, x_{2}, x_{3}\right)$ we want to look at

$$
\begin{equation*}
S^{\sharp} v(\lambda, x, y)=\psi\left(\lambda^{-\frac{2}{3}} y\right) \psi\left(\lambda^{-\frac{1}{3}} x\right) S v(\lambda, x, y) . \tag{7.25}
\end{equation*}
$$

The billiard ball maps for the model problem (7.21) have the following well known form

$$
\delta_{0}^{ \pm}\left(x_{1}, x_{3} ; \xi_{1}, \xi_{3}\right)=\left(x_{1} \pm 2\left(\frac{\xi_{1}}{\xi_{3}}\right)^{\frac{1}{2}}, x_{3} \mp \frac{2}{3}\left(\frac{\xi_{1}}{\xi_{3}}\right)^{\frac{3}{2}} ; \xi_{1}, \xi_{3}\right)
$$

In the new projective coordinates we obtain

$$
\gamma^{ \pm}(x, \xi)=\left(x \pm 2 \xi^{\frac{1}{2}}, \xi\right), \quad \xi \geq 0, \quad x=\xi_{3}^{\frac{1}{3}} x_{1}, \quad \xi=\xi_{3}^{-\frac{1}{3}} \xi_{1}
$$

that is $\gamma^{ \pm} \circ$ ( () $=$ ( ) $\circ \delta_{0}^{ \pm} \quad$ near $\Gamma_{0}$, where \%() $\left(x_{1}, x_{3} ; \xi_{1}, \xi_{3}\right)=\left(\xi_{3}^{\frac{1}{3}} x_{1}, \xi_{3}^{-\frac{1}{3}} \xi_{1}\right)$. Similarly, the chosen square root of the billiard ball map, $\mathcal{C}_{0}=\mathfrak{C}_{0}^{+} \cup \mathfrak{C}_{0}^{-}$,

$$
\begin{align*}
\mathfrak{e}_{0}^{ \pm}= & \left\{\left(\left(x_{1}, x_{3} ; \xi_{1}, \xi_{3}\right),\left(y_{1}, y_{3} ; \eta_{1}, \eta_{3}\right)\right): x_{1}-y_{1}= \pm\left(\frac{\xi_{1}}{\xi_{3}}\right)^{\frac{1}{2}}\right.  \tag{7.26}\\
& \left.\left(x_{3}-y_{3}\right)=\mp \frac{1}{3}\left(\frac{\xi_{1}}{\xi_{3}}\right)^{\frac{3}{2}}, \quad \xi_{i}=\eta_{i} \geq 0\right\}
\end{align*}
$$

can be rewritten as

$$
\mathscr{G}=\mathcal{G}^{+} \cup \mathcal{G}^{-}, \quad \mathcal{G}^{ \pm}=\left\{\left((x, \xi),\left(x^{\prime}, \xi^{\prime}\right)\right): \xi=\xi^{\prime} \geq 0, x-x^{\prime}= \pm \xi^{\frac{1}{2}}\right\}
$$

in the sense that $\mathscr{G}^{ \pm} \circ^{\%}(\mathcal{)}=\mathscr{(}) \circ \mathfrak{C}_{0}^{ \pm}$. We also define (near $\Gamma$ ):

$$
\text { (1) }+\left(x_{1}, x_{2}, x_{3} ; \xi_{1}, \xi_{2}, \xi_{3}\right)=\left(\xi_{3}^{\frac{1}{3}} x_{1}, \xi_{3}^{\frac{2}{3}} x_{2} ; \xi_{3}^{-\frac{1}{3}} \xi_{1}, \xi_{3}^{-\frac{2}{3}} \xi_{2}\right)
$$

and $\mathscr{H}$ such that

$$
\left.\mathscr{H}=\mathscr{H}^{+} \cup \mathscr{H}^{-}, \quad \mathscr{H}^{ \pm} \circ \mathscr{(}\right)=\mathscr{F}+\circ \mathscr{F}_{0}^{ \pm}
$$

with $\mathcal{F}_{0}^{ \pm} \subset T^{*} \mathbb{R}^{3} \backslash 0 \times T^{*} \mathbb{R}^{2}$ defined earlier and generated by $x_{1} \xi_{1}+x_{3} \xi_{3} \mp \frac{2}{3}(-\zeta(x, \xi))^{\frac{3}{2}}$. The relations $\mathscr{H}^{ \pm}$also take a very simple form

$$
\begin{equation*}
\mathscr{H}^{ \pm}=\left\{\left((x, y ; \xi, \eta) ;\left(x^{\prime}, \xi^{\prime}\right)\right): \xi=\xi^{\prime}, x=x^{\prime} \pm(\xi+y)^{\frac{1}{2}}, \eta=\mp(\xi+y)^{\frac{1}{2}}\right\} \tag{7.27}
\end{equation*}
$$

Let $\Lambda_{m}^{ \pm}=N^{*}\left\{x_{3}+\frac{1}{3 m^{2}} x_{1}^{3}=0\right\} \cap\left\{ \pm x_{1} \geq 0\right\} \subset T^{*} \mathbb{R}^{2} \backslash 0, m>00$, be the model boundary Lagrangians first introduced in (7.19):

$$
\begin{equation*}
\Lambda_{m}^{ \pm}=\left(\mathfrak{e}_{0}^{ \pm}\right)^{m} \Lambda_{0}, \quad \Lambda_{0}=N^{*}\left\{x_{1}=x_{3}=0\right\} \tag{7.28}
\end{equation*}
$$

We then consider $\Xi_{m}^{ \pm} \subset T^{*} \mathbb{R}$,

$$
\begin{equation*}
\Xi_{m}^{ \pm}=\vartheta\left(\Lambda_{m}^{ \pm}\right)=\left\{(x, \xi): \xi=m^{-2} x^{2}, \pm x \geq 0\right\}, m>0, \Xi_{0}=\{(0, \xi): \xi \in \mathbb{R}\} \tag{7.29}
\end{equation*}
$$

If $\widetilde{\Lambda}_{m}=\mathscr{F}_{0}\left(\Lambda_{m}\right), \Lambda_{m}=\Lambda_{m}^{+} \cup \Lambda_{m}^{-}$, then (7.27) immediately yields

$$
\begin{align*}
& \widetilde{\Xi}_{0}=\mathscr{\vartheta}+\left(\widetilde{\Lambda}_{0}\right)=\mathscr{H}\left(\Xi_{0}\right)=\left\{(x, y ; \xi, \eta): \eta^{2}-y-\xi=0, x+\eta=0\right\}  \tag{7.30}\\
& \widetilde{\Xi}_{m}=\vartheta\left(\widetilde{\Lambda}_{m}\right)=\mathscr{H}\left(\Xi_{m}\right)=\left\{(x, y ; \xi, \eta): \eta^{2}-y-\xi=0, m^{-2}(x+\eta)^{2}=\xi\right\}
\end{align*}
$$

when $m \neq 0$. Using this and the observation that, by homogeneity,

$$
\left(3 x_{3} \xi_{3}+2 x_{2} \xi_{2}+x_{1} \xi_{1}\right) \Gamma_{\tilde{\Lambda}_{m}}=0
$$

we obtain explicit expressions for $\widetilde{\Lambda}_{m}$ in terms of generators or as conormal bundles:

$$
\begin{aligned}
& \tilde{\Lambda}_{0}=N^{*}\left(\left\{x_{3}+\frac{1}{3} x_{1}^{3}-x_{1} x_{2}=0\right\}\right) \\
& \tilde{\Lambda}_{1}=\operatorname{cl}\left[N^{*}\left(\left\{x_{3} x_{1}+\frac{1}{12}\left(x_{2}-x_{1}^{2}\right)^{2}-\frac{1}{3} x_{1}\left(x_{2}+x_{1}^{2}\right)=0\right\} \backslash\left\{x_{1}=x_{2}=x_{3}=0\right\}\right)\right] \\
& \widetilde{\Lambda}_{2}= \\
& \quad \operatorname{cl}\left[N^{*}\left(\left\{27 x_{3}+9 x_{1} x_{2}+5 x_{1}^{3}\right)^{2}-16\left(3 x_{2}-x_{1}^{2}\right)^{3}=0\right\} \backslash\right. \\
& \left.\left.\quad\left\{3 x_{2}-x_{1}^{2}=27 x_{3}+9 x_{1} x_{2}+5 x_{1}^{3}=0\right\}\right)\right]
\end{aligned}
$$

The Lagrangians (7.28), $m \leq 3$, are grouped in two pairs related by the billiard ball map:

$$
\left\{\Lambda_{0}, \Lambda_{2}\right\}, \quad\left\{\Lambda_{1}, \Lambda_{3}\right\}, \quad \Lambda_{2}^{ \pm}=\delta_{0}^{ \pm} \Lambda_{0}, \quad \Lambda_{3}^{ \pm}=\delta_{0}^{ \pm} \Lambda_{1},
$$

and we study these separately as they cannot be simultaneous models - see (7.19) and the discussion following it.

The main result of this subsection is
Proposition 7.9. - The Dirichlet problem parametrix for (7.21), $T_{0}^{\sharp}$, given by (7.23) has the following mapping properties for $k$ even:

$$
\begin{aligned}
& T_{0}^{\sharp}: I_{k} L_{c}^{2}\left(\mathbb{R}^{2} ; \Lambda_{0}, J_{\Lambda_{0} \cap \Lambda_{2}}^{1} \Lambda_{2}, 2\right)+ I_{k} L_{c}^{2}\left(\mathbb{R}^{2} ; \Lambda_{2}^{+}, 2\right) \longrightarrow \\
& I_{k}^{b} L_{\mathrm{loc}}^{2}\left(\mathbb{R}_{+}^{3} ; \mathcal{F}_{0}\left(\Lambda_{1}\right), \mathfrak{F}_{0}\left(\Lambda_{1}\right) \cap \mathcal{F}_{0}\left(\Lambda_{0}\right)\right), \\
& T_{0}^{\sharp}: I_{k} L_{c}^{2}\left(\mathbb{R}^{2}, \Lambda_{1}^{-}, 2\right)+I_{k} L_{c}^{2}\left(\mathbb{R}^{2}, \Lambda_{3}^{+}, 2\right) \longrightarrow I_{k}^{b} L_{\mathrm{loc}}^{2}\left(\mathbb{R}_{+}^{3} ; \mathcal{F}_{0}\left(\Lambda_{2}\right), \mathcal{F}_{0}\left(\Lambda_{0}\right) \cap \mathcal{F}_{0}\left(\Lambda_{2}\right)\right),
\end{aligned}
$$

and

$$
T_{0}^{\sharp}: I_{k} L_{c}^{2}\left(\mathbb{R}^{2}, \Lambda_{1}^{+}, 2\right) \longrightarrow I_{k}^{b} L_{\mathrm{loc}}^{2}\left(\mathbb{R}_{+}^{3} ; \mathcal{F}_{0}\left(\Lambda_{0}\right) \sqcup \mathscr{F}_{0}\left(\Lambda_{\sharp}\right)\right),
$$

where $\Lambda_{\sharp} \subset T^{*} \mathbb{R}^{2} \backslash 0$ is any $C^{\infty}$ homogeneous Lagrangian tangent to $\Lambda_{1}$ at $\Lambda_{1} \cap \Lambda_{0}$.
We want to reduce the proof to an estimate for $S$ defined by (7.24) and for that we need a characterization of the spaces above. Thus we define

$$
\begin{aligned}
& I_{k} L^{2}\left(Z, \Xi_{m}^{ \pm}\right)=\left\{v \in L^{2}(Z):\left(\lambda D_{\lambda}\right)^{k_{0}}\left(D_{x}-m^{-2} x^{2}\right)^{k_{1}}\right. \\
& \left.\left(x\left(D_{x}-m^{-2} x^{2}\right)\right)^{k_{2}}(\chi(\mp x) x)^{k_{3}} v \in L^{2}, k_{0}+k_{2}+\frac{2}{3} k_{1}+\frac{1}{3} k_{3} \leq k\right\}
\end{aligned}
$$

for $k$ even, and by complex interpolation between the even indexed neighbours for $k$ odd. Similarly (but for all $k \in \mathbb{N}_{0}$ ),

$$
I_{k} L^{2}\left(Z, \Xi_{0}\right)=\left\{v \in L^{2}(Z):\left(\lambda D_{\lambda}\right)^{k_{0}}\left(x D_{x}\right)^{k_{1}} x^{k_{2}} v \in L^{2}, k_{0}+k_{1}+\frac{1}{3} k_{2} \leq k\right\}
$$

Definitions 4.3 and 4.5 easily give
Lemma 7.10. - If $W F(u) \subset \Gamma_{0}, u \in \mathcal{E}^{\prime}\left(\mathbb{R}^{2}\right)$ then

$$
\begin{aligned}
u & \in I_{k} L^{2}\left(\mathbb{R}^{2} ; \Lambda_{m}^{ \pm}, 2\right) \Longleftrightarrow W u \in I_{k} L_{\nu_{1 / 6}}^{2}\left(Z, \Xi_{m}^{ \pm}\right) \\
u & \in I_{k} L^{2}\left(\mathbb{R}^{2} ; \Lambda_{0}, J_{\Lambda_{0} \cap \Lambda_{1}}^{1} \Lambda_{2}, 2\right) \Longleftrightarrow W u \in I_{k} L_{\nu_{1 / 6}}^{2}\left(Z, \Xi_{0}\right),
\end{aligned}
$$

where $\nu_{\alpha}=\lambda^{-2 \alpha} d \lambda d x$.

We proceed similarly for $\widetilde{\Lambda}_{m}$ by defining

$$
\begin{gathered}
I_{k} L^{2}\left(Z_{+}, \widetilde{\Xi}_{0}\right)=\left\{u \in L^{2}\left(Z_{+}\right):\left(\lambda D_{\lambda}\right)^{k_{0}}\left(D_{y}^{2}-y-D_{x}\right)^{k_{1}}\left(D_{y}+x\right)^{k_{2}}\right. \\
\left.\quad\left(\left(x+D_{y}\right) D_{x}\right)^{k_{3}} u \in L^{2}\left(Z_{+}\right), k_{0}+\frac{2}{3} k_{1}+\frac{1}{3} k_{2}+k_{3} \leq k\right\}, \\
I_{k} L^{2}\left(Z_{+}, \widetilde{\Xi}_{m}\right)=\left\{u \in L^{2}\left(Z_{+}\right):\left(\lambda D_{\lambda}\right)^{k_{0}}\left(D_{y}^{2}-y-D_{x}\right)^{k_{1}}\left(D_{x}-m^{-2}\left(x+D_{y}\right)^{2}\right)^{k_{2}} .\right. \\
\left.\left[\left(x+D_{y}\right)\left(D_{x}-m^{-2}\left(x+D_{y}\right)^{2}\right)\right]^{k_{3}} \in L^{2}\left(Z_{+}\right), k_{0}+\frac{2}{3} k_{1}+\frac{2}{3} k_{2}+k_{3} \leq k\right\},
\end{gathered}
$$

for $k$ even and by complex interpolation for $k$ odd. The analogue of Lemma 7.10 does not hold in full generality but the following lemma is precisely what we need. For the notational convenience it is stated, as the rest of this subsection, for $\mathbb{R}_{+}^{3}$ but the generalization to $\mathbb{R}_{+}^{n}$ is easy.

Lemma 7.11. - If for some $\chi \in C_{0}^{\infty}\left(\mathbb{R}^{3}\right), \chi(0)=1, W F_{b}(\chi u) \subset \jmath_{*} \Gamma$ and $P_{0} u \in C^{\infty}\left(\overline{\mathbb{R}}_{+}^{3}\right)$, then for $k$ even

$$
\begin{array}{ll}
W_{+}(\varphi u) \in I_{k} L_{\nu_{\frac{1}{2}}^{+}}^{2}\left(Z, \widetilde{\Xi}_{m}\right) & \text { for any } \varphi \in C_{0}^{\infty}\left(\overline{\mathbb{R}}_{+}^{3}\right) \Longrightarrow \\
& u \in I_{k}^{b} L_{\mathrm{loc}}^{2}\left(\mathbb{R}_{+}^{3}, \mathcal{F}_{0}\left(\Lambda_{m}\right), \mathcal{F}_{0}\left(\Lambda_{m}\right) \cap \mathcal{F}_{0}\left(\Lambda_{0}\right)\right), \quad m \neq 0, \\
W_{+}(\varphi u) \in I_{k} L_{\nu_{\frac{1}{2}}^{+}}^{2}\left(Z, \widetilde{\Xi}_{0}\right) & \text { for any } \varphi \in C_{0}^{\infty}\left(\overline{\mathbb{R}}_{+}^{3}\right) \Longrightarrow \\
& u \in I_{k}^{b} L_{\mathrm{loc}}^{2}\left(\mathbb{R}_{+}^{3}, \mathscr{F}_{0}\left(\Lambda_{0}\right) \sqcup \mathcal{F}_{0}\left(\Lambda_{\sharp}\right)\right),
\end{array}
$$

$\nu_{\alpha}^{+}=\lambda^{-2 \alpha} d \lambda d x d y$, and where $\Lambda_{\sharp} \subset T^{*} \mathbb{R}^{2} \backslash 0$ is any $C^{\infty}$ homogeneous Lagrangian tangent to $\Lambda_{1}$ at $\Lambda_{0} \cap \Lambda_{1}$.

Proof. - We start with the case $m \neq 0$. Near $\jmath_{*} \Gamma,{ }^{b} \mathscr{\Re}_{1}^{\sharp}\left(\mathscr{F}_{0}\left(\Lambda_{m}\right), \mathfrak{F}_{0}\left(\Lambda_{m}\right) \cap \mathfrak{F}\left(\Lambda_{0}\right)\right)$ is generated by (see (7.29))

$$
\begin{align*}
& p_{0}(x, \xi)=\xi_{2}^{2}-x_{2} \xi_{3}^{2}-\xi_{3} \xi_{1}, \quad a_{1}(x, \xi)=\sum_{j=1}^{3} j x_{j} \xi_{j}  \tag{7.31}\\
& a_{2}(x, \xi)^{2}, \quad a_{3}(x, \xi)=\left(x_{1} \xi_{3}+\xi_{2}\right) a_{2}(x, \xi)
\end{align*}
$$

where $a_{2}(x, \xi)=\xi_{1} \xi_{3}-m^{-2}\left(x_{1} \xi_{3}+\xi_{2}\right)^{2}$. We shall denote by $P_{0}$ and $A_{i}$ the differential operators corresponding to $p_{0}, a_{i}$. Since $P_{0} u \equiv 0$, it suffices to have

$$
\begin{equation*}
A_{1}^{\alpha_{1}} A_{2}^{\alpha_{2}} A_{3}^{\alpha_{3}} u \in \bar{H}_{\left(-\frac{3}{2} \alpha_{2}-2 \alpha_{3}\right)}^{\mathrm{loc}}\left(\mathbb{R}_{+}^{3}\right), \quad \alpha_{1}+\frac{1}{2} \alpha_{2}+\alpha_{3} \leq k=2 l \tag{7.32}
\end{equation*}
$$

In fact, $\left[P_{0}, A_{j}\right]=-4 i \delta_{1 j} P_{0}$ so that $P_{0} A_{1}^{\alpha_{1}} A_{2}^{\alpha_{2}} A_{3}^{\alpha_{3}} u \equiv 0$. Thus, (7.32) implies that $A_{1}^{\alpha_{1}} A_{2}^{\alpha_{2}} A_{3}^{\alpha_{3}} u \in \bar{H}_{\left(\alpha_{1}+2 \alpha_{2}+3 \alpha_{3},-\alpha_{1}-\frac{7}{2} \alpha_{2}-5 \alpha_{3}\right)}^{\text {loc }}\left(\mathbb{R}_{+}^{3}\right)$ which in turn shows that $B A_{1}^{\alpha_{1}} A_{2}^{\alpha_{2}} A_{3}^{\alpha_{3}} u \in L_{\text {loc }}^{2}\left(\mathbb{R}_{+}^{3}\right)$ for $B \in \Psi_{b}^{-\frac{3}{2} \alpha_{2}-2 \alpha_{3}, \alpha_{1}+2 \alpha_{1}+3 \alpha_{3}}\left(\mathbb{R}_{+}^{3}\right)$ and gives the defining condition for $I_{k}^{b} L_{\text {loc }}^{2}\left(\mathbb{R}_{+}^{3}, \mathscr{F}_{0}\left(\Lambda_{m}\right), \mathscr{F}_{0}\left(\Lambda_{m}\right) \cap \mathscr{F}_{0}\left(\Lambda_{0}\right)\right)$ :

$$
\left(B_{1} A_{1}\right)^{\alpha_{1}}\left(B_{2} A_{2}\right)^{\alpha_{2}}\left(B_{3} A_{3}\right)^{\alpha_{3}} u \in L_{\mathrm{loc}}^{2}\left(\mathbb{R}_{+}^{3}\right), \quad \alpha_{1}+\frac{1}{2} \alpha_{2}+\alpha_{3} \leq k
$$

where $B_{1} \in \Psi_{b}^{0}\left(\mathbb{R}_{+}^{3}\right), B_{2} \in \Psi_{b}^{-\frac{3}{2}}\left(\mathbb{R}_{+}^{3}\right), B_{3} \in \Psi_{b}^{-2}\left(\mathbb{R}_{+}^{3}\right)$. On the other hand the definition of $I_{k} L_{\nu_{\frac{1}{2}}^{+}}^{2}\left(Z_{+}, \widetilde{\Xi}_{m}\right)$ and the assumptions on $u$ imply that

$$
\begin{equation*}
A_{1}^{\beta_{1}} A_{2}^{\beta_{2}} A_{3}^{\beta_{3}} u \in \bar{H}_{\left(0,-\frac{4}{3} \beta_{2}-2 \beta_{3}\right)}^{\text {loc }}\left(\mathbb{R}_{+}^{3}\right), \quad \beta_{1}+\frac{2}{3} \beta_{2}+\beta_{3} \leq k=2 l . \tag{7.33}
\end{equation*}
$$

Let us observe that $\alpha_{2}$ in (7.32) can be assumed to be even and thus we have two cases

$$
\alpha_{2}=4 m+p, \quad p=0 \quad \text { or } \quad p=2 .
$$

If $p=0$ then

$$
A_{1}^{\alpha_{1}} A_{2}^{\alpha_{2}} A_{3}^{\alpha_{3}} u=A_{1}^{\alpha_{1}} A_{2}^{3 m} A_{2}^{m} A_{3}^{\alpha_{3}} u \in \bar{H}_{\left(-2 m,-4 m-2 \alpha_{3}\right)}^{\mathrm{loc}}\left(\mathbb{R}_{+}^{3}\right) \subset \bar{H}_{\left(-\frac{3}{2} \alpha_{2}-2 \alpha_{3}\right)}^{\mathrm{loc}}\left(\mathbb{R}_{+}^{3}\right)
$$

by (7.33) with $\beta_{i}=\alpha_{i}, i=1,3$ and $\beta_{2}=3 \mathrm{~m}$. Thus it remains to analyse the case $p=2$ which is more involved and which by the above ( $p=0$ ) argument reduces to $\alpha_{2}=2, k=2$ in (7.32). Let us note the following identities:

$$
\begin{equation*}
A_{2}=D_{x_{1}} D_{x_{3}}-m^{-2} A_{0}^{2}, A_{3}=A_{0} A_{2}, a_{0}(x, \xi)=x_{1} \xi_{3}+\xi_{2}, a_{0}=\sigma_{1}\left(A_{0}\right) \tag{7.34}
\end{equation*}
$$

The desired property (7.32) reduces to

$$
\begin{equation*}
A_{i} A_{2}^{2} u \in \bar{H}_{\left(-3-2 \delta_{3 i}\right)}^{\mathrm{loc}}\left(\mathbb{R}_{+}^{3}\right), \quad i=1,3 \tag{7.35}
\end{equation*}
$$

which for $i=3$ follows by writing $A_{3} A_{2}^{2} u=A_{0} A_{2}^{3} u$ with $A_{2}^{3} u \in \bar{H}_{(0,-4)}^{\text {loc }}$ by (7.33). Thus in view of the assumption $W F_{b}(u) \subset \jmath_{*} \Gamma$, it remains to establish

$$
\begin{equation*}
\left\|\left\langle D_{x_{3}}\right\rangle^{-3} \chi_{0} A_{1} A_{2}^{2} u\right\|_{L^{2}\left(\mathbb{R}_{+}^{3}\right)}^{2} \leq C \sum_{\beta_{1}+\frac{2}{3} \beta_{2}+\beta_{3} \leq 2}\left\|\left\langle D_{x_{3}}\right\rangle^{-\frac{4}{3} \beta_{2}-2 \beta_{3}} \chi_{1} A_{1}^{\beta_{1}} A_{2}^{\beta_{2}} A_{3}^{\beta_{3}} u\right\|_{L^{2}\left(\mathbb{R}_{+}^{3}\right)}^{2} \tag{7.36}
\end{equation*}
$$

where $\chi_{i}(x, D) \in \Psi^{0}\left(\partial \mathbb{R}_{+}^{3}\right), \operatorname{supp} \chi_{i}, i=0,1$ is in a neighbourhood of $\Gamma_{0}$ (which is a conic neighbourhood of $(0 ;(0,0,1))$ and $\chi_{1} \equiv 1$ on $\operatorname{supp} \chi_{0}$. We will denote the right hand side in (7.36) by $M(u)$. We first obtain an a priori inequality in which we assume $u \in C^{\infty}\left(\mathbb{R}_{+}^{3}\right)$ and start by using (7.34) to rewrite the left hand side of (7.36) as

$$
\left\langle D_{x_{3}}\right\rangle^{-3}\left\langle\chi_{0} A_{1} A_{2}^{2} u,\left\langle D_{x_{3}}\right\rangle^{-3} \chi_{0} A_{1} A_{2} D_{x_{1}} D_{x_{3}} u-m^{-2}\left\langle D_{x_{3}}\right\rangle^{-3} \chi_{0} A_{1} A_{2} A_{0}^{2} u\right\rangle
$$

which modulo commutator terms bounded by $M(B u), B \in \Psi_{b}^{0, K}\left(\mathbb{R}_{+}^{3}\right)$ for some $K$ (see the proof of Proposition 4.5), is equal to

$$
\begin{aligned}
& \left\langle\left\langle D_{x_{3}}\right\rangle^{-3} \chi_{0} D_{x_{1}} D_{x_{3}} A_{1} A_{2}^{2} u,\left\langle D_{x_{3}}\right\rangle^{-3} \chi_{0} A_{1} A_{2} u\right\rangle- \\
& m^{-2}\left\langle\left\langle D_{x_{3}}\right\rangle^{-3} \chi_{0} A_{1} A_{2}^{2} u,\left\langle D_{x_{3}}\right\rangle^{-3} \chi_{0} A_{1} A_{3} A_{0} u\right\rangle \equiv \\
& \left\langle\left\langle D_{x_{3}}\right\rangle^{-3} \chi_{0} A_{2}^{3} u,\left\langle D_{x_{3}}\right\rangle^{-3} \chi_{0} A_{1}^{2} A_{2} u\right\rangle+ \\
& m^{-2}\left\langle\left\langle D_{x_{3}}\right\rangle^{-3} \chi_{0} A_{0} A_{3} A_{1} A_{2} u,\left\langle D_{x_{3}}\right\rangle^{-3} \chi_{0} A_{1} A_{2} u\right\rangle- \\
& m^{-2}\left\langle\left\langle D_{x_{3}}\right\rangle^{-3} \chi_{0} A_{1} A_{2}^{2} u,\left\langle D_{x_{3}}\right\rangle^{-3} \chi_{0} A_{1} A_{3} A_{0} u\right\rangle .
\end{aligned}
$$

The integrations by parts without boundary contributions were allowed as we used only the tangential derivatives and $x_{2} D_{x_{2}}$ in $A_{1}$ which vanishes at the boundary. Thus the following estimate holds

$$
\begin{aligned}
\left\|\left\langle D_{x_{3}}\right\rangle^{-3} \chi_{0} A_{1} A_{2}^{2} u\right\| \leq & \sum_{j=1}^{J} M\left(B_{j} u\right)+\left\|\left\langle D_{x_{3}}\right\rangle^{-4} \chi_{0} A_{2}^{3} u\right\|\left\|\left\langle D_{x_{3}}\right\rangle^{-2} \chi_{0} A_{1}^{2} A_{2} u\right\|+ \\
& \left\|\left\langle D_{x_{3}}\right\rangle^{-4} \chi_{0} A_{0} A_{3} A_{1} A_{2} u\right\|\left\|\left\langle D_{x_{3}}\right\rangle^{-2} \chi_{0} A_{1} A_{2} u\right\|+ \\
& \left\|\left\langle D_{x_{3}}\right\rangle^{-3} \chi_{0} A_{1} A_{2}^{2} u\right\|\left\|\left\langle D_{x_{3}}\right\rangle^{-3} \chi_{0} A_{1} A_{3} A_{0} u\right\|,
\end{aligned}
$$

where the norms are in $L^{2}\left(\mathbb{R}_{+}^{3}\right)$. An approximation argument gives

$$
\left\|\left\langle D_{x_{3}}\right\rangle^{-3} \chi_{0} A_{1} A_{2}^{2} u\right\| \leq 2 \sum_{j=1}^{\tilde{J}} M\left(B_{j} u\right)
$$

for any $u$ satisfying (7.33) with $k=2$. As in the proof of sufficiency of (7.32) $P_{0} u \equiv 0$ implies now that $M(B u) \leq C M(u)$ for $B \in \Psi_{b}^{0, K}\left(\mathbb{R}_{+}^{3}\right)$. In fact, commuting $B$ through gives

$$
M(B u) \leq C \quad \sum_{\beta_{1}+\frac{2}{3} \beta_{2}+\beta_{3} \leq 2} \sum_{j=1}^{J_{\beta}}\left\|B_{j}^{\beta} A_{1}^{\beta_{1}} A_{2}^{\beta_{2}} A_{3}^{\beta_{3}} u\right\|^{2}
$$

where $B_{j}^{\beta} \in \Psi_{b}^{-\frac{4}{3} \beta_{2}-3 \beta_{3}, K+\beta_{1}+2 \beta_{2}+3 \beta_{3}}\left(\mathbb{R}_{+}^{3}\right)$ and hence we need to show that for $E \in \Psi_{b}^{m, k}\left(\mathbb{R}_{+}^{3}\right)$ and $v$ such that $P_{0} v \in C^{\infty}$ near the boundary,

$$
\|E v\| \leq C\|v\|_{(0, m)} \quad \text { for } \quad m \leq 0 .
$$

This in turn follows from $\|E v\| \leq C\|v\|_{(k, m-k)}, m-k \leq 0$, and Theorem B.2.9 of [14]. Thus we obtain (7.35), $i=1$, which concludes the proof of (7.32) and of the lemma for $m>0$.
For the case $m=0$ we need to describe the generators of ${ }^{b} 9 \mathbb{N}_{1}^{\sharp}\left(\mathscr{F}_{0}\left(\Lambda_{0}\right) \sqcup \mathscr{F}_{0}\left(\Lambda_{\sharp}\right)\right)$ near $\jmath_{*} \Gamma$. The Lagrangian $\mathscr{F}\left(\Lambda_{0}\right)=\widetilde{\Lambda}_{0}$ is given near $\Gamma$, a conic neighbourhood of $(0 ;(0,0,1))$ by the zeros of $p_{0}, a_{0}, a_{1}$ defined in (7.31) and (7.34). Since $\Lambda_{\sharp}$ is tangent to $\Lambda_{1}$ at $\Lambda_{0} \cap \Lambda_{1}$, it follows that $\Lambda_{\sharp}=N^{*}\left\{x_{3}+x_{1}^{3} f\left(x_{1}\right)=0\right\}, f \in C^{\infty}$. Thus, a computation based on this and the definition of $\mathscr{F}_{0}$ shows that, near $\Gamma, \mathscr{F}_{0}\left(\Lambda_{\sharp}\right)$ is given by the zeros of $p_{0}, a_{1}^{\sharp}, a_{2}^{\sharp}$ where

$$
\begin{array}{r}
a_{1}^{\sharp}(x, \xi)=a_{1}(x, \xi)+\xi_{3}^{-2} a_{0}(x, \xi)^{3} h_{1}(x, \xi),  \tag{7.37}\\
a_{2}^{\sharp}(x, \xi)=\xi_{1} \xi_{3}-a_{0}(x, \xi)^{2} h_{2}(x, \xi),
\end{array}
$$

with

$$
\begin{aligned}
& h_{1}(x, \xi)=\xi_{3}^{-1} a_{0}(x, \xi) f^{\prime}\left(x_{1}+\xi_{3}^{-1} \xi_{2}\right) / 3 \\
& h_{2}(x, \xi)=f\left(x_{1}+\xi_{3}^{-1} \xi_{2}\right)+\xi_{3}^{-1} a_{0}(x, \xi) f^{\prime}\left(x_{1}+\xi_{3}^{-1} \xi_{2}\right)
\end{aligned}
$$

To obtain the generators of ${ }^{b} 9 \Re_{1}^{\sharp}\left(\mathcal{F}_{0}\left(\Lambda_{\sharp}\right)\right)$ near $J_{*} \Gamma$ we write

$$
h_{i}(x, \xi)=h_{i}^{e}\left(x, \xi_{1}, \xi_{3}\left(\xi_{2} / \xi_{3}\right)^{2}, \xi_{3}\right)+\left(\xi_{2} / \xi_{3}\right) h_{i}^{o}\left(x, \xi_{1}, \xi_{3}\left(\xi_{2} / \xi_{3}\right)^{2}, \xi_{3}\right)
$$

so that using $p_{0}, h_{i}$ in $a_{i}^{\sharp}$ can be replaced by

$$
h_{i}^{\sharp}(x, \xi)=h_{i}^{e}\left(x ; \xi_{1}, \xi_{3} x_{1}+\xi_{1}, \xi_{3}\right)+\xi_{2} \xi_{3}^{-1} h_{i}^{o}\left(x ; \xi_{1}, \xi_{3} x_{1}+\xi_{1}, \xi_{3}\right)
$$

which near $\jmath_{*} \Gamma$ are in $S_{\mathrm{hg}}^{0,1}\left({ }^{b} T^{*} \overline{\mathbb{R}}_{+}^{3}\right)$, so that the corresponding $a_{1}^{\sharp} \in S_{\mathrm{hg}}^{1,4}\left({ }^{b} T^{*} \overline{\mathbb{R}}_{+}^{3}\right)$ and $a_{2}^{\sharp} \in S_{\mathrm{hg}}^{2,3}\left({ }^{b} T^{*} \overline{\mathbb{R}}_{+}^{3}\right)$.
We observe that we can find operators $H_{i} \in \Psi_{b}^{0,1}\left(\mathbb{R}_{+}^{3}\right)$ such that ${ }^{b} \sigma_{0,1}\left(H_{i}\right)=h_{i}^{\sharp}$ in $\jmath \Gamma$ and

$$
\begin{equation*}
\left[P_{0}, H_{i}\right]=D_{x_{2}}^{2} B_{-1}^{1}+B_{-1}^{2} D_{x_{2}}^{2}+B_{0}^{2} D_{x_{2}}+B_{1}^{1}, \quad B_{i}^{j} \in \Psi_{b}^{i}\left(\mathbb{R}_{+}^{3}\right), W F_{b}\left(B_{i}^{j}\right) \cap \jmath \Gamma=\emptyset \tag{7.38}
\end{equation*}
$$

where $W F_{b}(B)$ denotes the essential support of $B$ as a $b$-pseudo-differential operator. The wave front set conditions on $u$ and $B$ imply that $B u \in \dot{\mathscr{Q}}\left(\mathbb{R}_{+}^{3}\right)$. Since $u \in \mathscr{G} \tau\left(\mathbb{R}_{+}^{3}\right)$ it then follows that $B u \in C^{\infty}\left(\overline{\mathbb{R}}_{+}^{3}\right)$ as $B: \mathfrak{\tau}\left(\mathbb{R}_{+}^{3}\right) \subset \mathfrak{Q}^{\prime}\left(\mathbb{R}_{+}^{3}\right) \rightarrow \mathfrak{Q}^{\prime}\left(\mathbb{R}_{+}^{3}\right)$ and $\mathfrak{Q}^{\prime}\left(\mathbb{R}_{+}^{3}\right) \cap \dot{\mathfrak{d}}\left(\mathbb{R}_{+}^{3}\right)=C^{\infty}\left(\overline{\mathbb{R}}_{+}^{3}\right)$ (see the references given above and also Subsection 7.3 below). Hence $\left[P_{0}, H_{i}\right] u \in C^{\infty}\left(\overline{\mathbb{R}}_{+}^{3}\right)$ and we can commute $P_{0}$ through as in the proof of (7.32).
From (7.37) we obtain the generators of ${ }^{b} \mathscr{N}_{1}^{\sharp}\left(\mathscr{F}_{0}\left(\Lambda_{0}\right) \sqcup \mathcal{F}_{0}\left(\Lambda_{\sharp}\right)\right)$ in a neighbourhood of $\jmath_{*} \Gamma$ :

$$
\xi_{3}^{-1} p_{0}, \quad a_{1}^{\sharp}, \quad a_{2}^{\sharp} a_{0} .
$$

The assumption $W_{+} u \in I_{k} L_{\nu_{\frac{1}{2}}}^{2}\left(Z_{+}, \widetilde{\Xi}_{0}\right)$ implies

$$
\begin{equation*}
A_{0}^{\beta_{0}} A_{1}^{\beta_{1}}\left(D_{x_{1}} A_{0}\right)^{\beta_{2}} u \in \bar{H}_{\left(0,-\frac{2}{3} \beta_{0}-\beta_{2}\right)}^{\mathrm{loc}}, \quad \frac{1}{3} \beta_{0}+\beta_{1}+\beta_{2} \leq k \tag{7.39}
\end{equation*}
$$

Thus, as in the case of (7.32), it suffices to prove that (7.39) implies

$$
\begin{equation*}
\left(A_{1}^{\sharp}\right)^{\alpha_{1}}\left(A_{2}^{\sharp} A_{0}\right)^{\alpha_{2}} u \in \bar{H}_{\left(-m,-2 \alpha_{2}+m\right)}^{\mathrm{loc}}, \quad{ }^{b} \sigma_{1,4}\left(A_{1}^{\sharp}\right)=a_{1}^{\sharp},{ }^{b} \sigma_{2,3}\left(A_{2}^{\sharp}\right)=a_{2}^{\sharp}, \tag{7.40}
\end{equation*}
$$

for some $m$. However, (7.37) shows that

$$
A_{1}^{\sharp} \equiv\left\langle D_{x_{3}}\right\rangle^{-2}\left(D_{x_{3}}^{2} A_{1}+A_{0}^{3} H_{1}\right), \quad A_{2}^{\sharp} A_{0} \equiv D_{x_{3}} D_{x_{1}} A_{0}-A_{0}^{3} H_{2}, \quad H_{i} \in \Psi_{b}^{0,1}\left(\overline{\mathbb{R}}_{+}^{3}\right) .
$$

and thus (7.40) follows from (7.39).
The proof of Proposition 7.9 reduces to the proof of
Proposition 7.12. - If $S^{\sharp}$ is defined by (7.25) then for $k$ even

$$
\begin{equation*}
S^{\sharp}: I_{k} L_{\nu_{\alpha}}^{2}\left(Z ; \Xi_{m}^{ \pm}\right) \longrightarrow I_{k} L_{\nu_{\alpha+\frac{1}{3}}^{+}}^{2}\left(Z_{+} ; \widetilde{\Xi}_{m_{ \pm 1}}\right) \tag{7.41}
\end{equation*}
$$

$$
\begin{equation*}
S^{\sharp}: I_{k} L_{\nu_{\alpha}}^{2}\left(Z ; \Xi_{0}\right) \longrightarrow I_{k} L_{\nu_{\alpha+\frac{1}{3}}^{+}}^{2}\left(Z_{+} ; \widetilde{\Xi}_{1}\right) \tag{7.42}
\end{equation*}
$$

where $\nu_{\alpha}=\lambda^{-2 \alpha} d \lambda d x, \nu_{\beta}^{+}=\lambda^{-2 \beta} d \lambda d x d y$.
We start with the second mapping property where we can give a direct proof.
Proof. - Proof of (7.42) Since $\lambda$ does not appear in $S$ other than in the cut-off function $\psi$, the stability under $\lambda D_{\lambda}$ is clear and for simplicity of notation we shall omit that variable. Thus we consider

$$
\widetilde{S} \widehat{v}(\xi, y)=\int S v(x, y) e^{-i x \xi} d x, \quad \widetilde{S} \widehat{v}(\xi, y)=\psi\left(\lambda^{-\frac{2}{3}} \xi\right) \frac{A_{+}(-\xi-y)}{A_{+}(-\xi)} \widehat{v}(\xi)
$$

and recalling the definition of $I_{k} L_{\nu_{+}}^{2}\left(Z_{+}, \widetilde{\Xi}_{1}\right)$ we want to apply to it the operators

$$
\xi-\left(D_{y}-D_{\xi}\right)^{2}, \quad\left(\xi-\left(D_{y}-D_{\xi}\right)^{2}\right)\left(D_{y}-D_{\xi}\right)
$$

with weights $\frac{2}{3}$ and 1 , respectively, while the stability under $D_{y}^{2}-y-\xi$ is clear from the Airy equation. We claim that

$$
\left(\left(\xi-\left(D_{y}-D_{\xi}\right)^{2}\right)\left(D_{y}-D_{\xi}\right)\right)^{k_{2}}\left(\xi-\left(D_{y}-D_{\xi}\right)^{2}\right)^{k_{1}} \widetilde{S} \widehat{v}(\xi, y) \equiv \widetilde{S} \widehat{v}_{k_{1}, k_{2}}(\xi, y)
$$

where $v_{k_{1}, k_{2}} \in L_{\nu}^{2}(Z)$ if $v \in I_{k} L_{\nu}^{2}\left(Z, \Xi_{0}\right)$ and $\frac{2}{3} k_{1}+k_{2} \leq k$ and where we omitted the irrelevant terms with differentiation falling on $\psi$. In fact, we can proceed by induction, noting that the order in the iteration does not matter: $\left(\xi-\left(D_{y}-D_{\xi}\right)^{2}\right) \widetilde{S} \widehat{v}_{k_{1}-1,0}(\xi, y) \equiv$ $\widetilde{S} \widehat{v}_{k_{1} 0}(\xi, y)$, where by a simple computation

$$
\begin{equation*}
\widehat{v}_{k_{1} 0}(\xi)=2\left(D \Phi_{+}\right)(-\xi) \widehat{v}_{k_{1}-1,0}(\xi)-2 \Phi_{+}(-\xi) D_{\xi} \widehat{v}_{k_{1}-1,0}(\xi)-D_{\xi}^{2} \widehat{v}_{k_{1}-1,0}(\xi) \tag{7.43}
\end{equation*}
$$

Here

$$
\Phi_{+}(t)=\frac{1}{i} \frac{A_{+}^{\prime}(t)}{A_{+}(t)}, \quad-D_{t} \Phi_{+}(t)=\Phi_{+}(t)^{2}+t, \quad \Phi_{+} \in S^{\frac{1}{2}}(\mathbb{R})
$$

Thus, to obtain the boundedness of the second term in (7.43), we need the stability of $\widehat{v}_{k_{1}-1,0}(\xi), \frac{2}{3} k_{1} \leq k$ under $\langle\xi\rangle^{\frac{1}{2}} D_{\xi}$ with weight $\frac{2}{3}$. This follows easily from the stability under $\xi D_{\xi}$ with weight $1, D_{\xi}$ with weight $1 / 3$ and an interpolation argument (see the proof of Lemma 7.13).

We now turn to $\left(D_{y}-D_{\xi}\right)\left(\xi-\left(D_{y}-D_{\xi}\right)^{2}\right) \widetilde{S} \widehat{v}_{k_{1}, k_{2}-1}(\xi, y) \equiv \widetilde{S} \widehat{v}_{k_{1} k_{2}}(\xi, y)$ where now

$$
\begin{aligned}
\widehat{v}_{k_{1} k_{2}}(\xi)= & -2 D^{2} \Phi_{+}(-\xi) \widehat{v}_{k_{1}, k_{2}-1}(\xi)+4 D \Phi_{+}(-\xi) D_{\xi} \widehat{v}_{k_{1}, k_{2}-1}(\xi) \\
& -2 \Phi_{+}(-\xi) D_{\xi}^{2} \widehat{v}_{k_{1}, k_{2}-1}(\xi)-D_{\xi}^{3} \widehat{v}_{k_{1}, k_{2}-1}(\xi)
\end{aligned}
$$

and we use the stability of $\widehat{v}$ under $\langle\xi\rangle^{\frac{1}{2}} D_{\xi}^{2}$ with weight 1 . The proof is concluded by observing that $\psi\left(\lambda^{-\frac{2}{3}} y\right) \widetilde{S}=\mathcal{O}(1): \lambda^{\frac{1}{6}} L^{2}\left(\mathbb{R}_{\xi}\right) \longrightarrow \lambda^{\frac{1}{2}} L^{2}\left(\mathbb{R}_{\xi} \times \mathbb{R}_{+y}\right)$.

The presence of Lagrangians with boundaries in (7.41) presents difficulties which seem to prevent a computational proof similar to the one above, and the following characterization of the partial Fourier transform of $I_{k} L^{2}\left(Z, \Xi_{m}^{ \pm}\right)$will allow us to overcome them. We need yet another two spaces of functions:

$$
S_{k}^{\alpha} L^{2}(Z)=\left\{a \in \lambda^{\alpha} L^{2}(Z):\left(\lambda D_{\lambda}\right)^{k_{0}}\left(\xi D_{\xi}\right)^{k_{1}} D_{\xi}^{k_{2}} a(\lambda, \xi) \in \lambda^{\alpha} L^{2}(Z), \quad k_{0}+k_{1}+\frac{1}{3} k_{2} \leq k\right\}
$$

and

$$
\mathfrak{s}_{k}^{\alpha} L^{2}(Z)=\left\{a \in \lambda^{\alpha} L^{2}(Z):\left(\lambda D_{\lambda}\right)^{k_{0}} \xi^{k_{1}} D_{\xi}^{k_{2}} a \in \lambda^{\alpha} L^{2}(Z), \quad k_{0}+\frac{2}{3} k_{1}+\frac{1}{3} k_{1} \leq k\right\}
$$

with the latter space defined only for even $k$ and then for odd $k$ by complex interpolation between $k-1$ and $k+1$.

Using this we now have the crucial
Lemma 7.13. - Let $\widehat{v}(\lambda, \xi)$ be the Fourier transform of $v(\lambda, x)$ in the second variable. Then for $k$ even

$$
\begin{equation*}
v \in I_{k} L_{\nu_{\alpha}}^{2}\left(Z ; \Xi_{m}^{ \pm}\right) \Longleftrightarrow \widehat{v}(\lambda, \xi)=g(\lambda, \xi)+e^{\mp \frac{2}{3} i m \xi^{\frac{3}{2}}} f(\lambda, \xi) \tag{7.44}
\end{equation*}
$$

where $f \in S_{k}^{\alpha} L^{2}(Z)$, supp $f \subset\{\xi \geq 1\}$ and

$$
g \in \begin{cases}S_{k}^{\alpha} L^{2}(Z) & m=0 \\ \mathcal{S}_{k}^{\alpha} L^{2}(Z) & m>0\end{cases}
$$

$\Xi_{0}^{ \pm} \equiv \Xi_{0}, \quad \nu_{\alpha}=\lambda^{-2 \alpha} d \lambda d x$.
Proof. - We observe that the case $m=0$ follows immediately from the definition and we shall first prove that the left hand side in (7.44) implies the right hand side, that is, we assume that $v \in I_{k} L_{\nu_{\alpha}}^{2}\left(Z ; \Xi_{m}^{ \pm}\right)$for $k=2 l$ and $m>0$. To simplify the notation we will allow $m \in \mathbb{Z}$ and define $\Xi_{m} \stackrel{\text { def }}{=} \Xi_{|m|}^{\operatorname{sgn}(m)}$, and also put $\alpha=\frac{1}{6}$. Thus we take $v \in I_{2 l} L_{\nu_{\frac{1}{6}}}^{2}(Z, \Xi$,$) .$
Let $\chi \in C^{\infty}(\mathbb{R})$, supp $\chi \subset[1, \infty)$ be such that $\chi \equiv 1$ for $x>2$ and let us define

$$
v_{+}(x)=\chi(\operatorname{sgn}(m) x) v(x), \quad v_{-}=v-v_{+}
$$

The term $v_{-}$satisfies the estimates $\left(\lambda D_{\lambda}\right)^{k_{0}} x^{k_{1}} D_{x}^{k_{2}} v_{-} \in \lambda^{1 / 6} L^{2}(Z)$ for $k_{0}+\frac{1}{3} k_{1}+$ $\frac{2}{3} k_{2} \leq 2 l$, which immediately implies that $\hat{v}_{-} \in \mathscr{S}_{2 l}^{1 / 6} L^{2}(Z)$.
We can easily construct a canonical transformation $\chi_{m}: T^{*} \mathbb{R}^{2} \backslash 0 \longrightarrow T^{*} \mathbb{R}^{2} \backslash 0$ such that $\chi_{m}: N^{*}\left\{x_{1}=0\right\} \longrightarrow \Lambda_{m}$ :

$$
\begin{equation*}
\chi_{m}:\left(x_{1}, x_{2} ; \xi_{1}, \xi_{2}\right) \longmapsto\left(x_{1}-\frac{1}{3 m^{2}} x_{2}^{3}, x_{2} ; \xi_{1}, \xi_{2}+\frac{1}{m^{2}} x_{2}^{2} \xi_{1}\right)=\left(y_{1}, y_{2} ; \eta_{1}, \eta_{2}\right) \tag{7.45}
\end{equation*}
$$

which is generated by $\phi_{m}=\frac{1}{3 m^{2}} y_{2}^{3} \xi_{1}+y_{2} \xi_{2}+y_{1} \xi_{1}$. The definition of $I_{k} L_{\nu_{\alpha}}^{2}\left(Z ; \Xi_{m}^{ \pm}\right)$ then shows that (see Lemma 7.10)

$$
v(\lambda, x)=e^{\frac{1}{3} i m^{-2} x^{3}} W u_{1}(\lambda, x), \quad u_{1} \in I_{k} L^{2}\left(\mathbb{R}^{2} ; N^{*}\left\{x_{1}=0\right\} \cap\left\{ \pm x_{2} \geq 0\right\}, 2\right)
$$

Consequently, if $w(\lambda, x)=e^{-\frac{1}{3 m^{2}} i x^{3}} v_{+}(\lambda, x)$ then since $\hat{v}_{-} \in \mathcal{S}_{2 l}^{1 / 6} L^{2}(Z)$,

$$
\begin{equation*}
\left(\lambda D_{\lambda}\right)^{k_{0}}\left(x D_{x}\right)^{k_{1}} D_{x}^{k_{2}} w \in \lambda^{\frac{1}{6}} L^{2}(Z), \quad k_{0}+k_{1}+\frac{2}{3} k_{2} \leq 2 l, \operatorname{supp} w \subset \operatorname{sgn}(m)[1, \infty) \tag{7.46}
\end{equation*}
$$

Strictly speaking we should now replace $w \in \lambda^{1 / 6} L^{2}(Z)$ by a sequence $w_{j} \in \lambda^{1 / 6} L^{2}([1, \infty) ; \varsigma(\mathbb{R})), w_{j} \rightarrow w$ with the estimate (7.46) satisfied uniformly, but for simplicity we shall write $w$ everywhere. With this understanding we have

$$
\begin{equation*}
\widehat{v}_{+}(\lambda, \xi)=\int e^{\frac{1}{3} i m^{-2} x^{3}-i x \xi} w(\lambda, x) d x \tag{7.47}
\end{equation*}
$$

and we first consider $\xi>1$ where we define $f(\lambda, \xi)=\chi(\xi) \exp \left(\frac{2}{3} i m \xi^{\frac{3}{2}}\right) \widehat{v}_{+}(\lambda, \xi)$. We will check the stability under $D_{\xi}$ and $\xi D_{\xi}$ ( $\lambda D_{\lambda}$ is clear) by induction for more general $f$ 's of the form

$$
\begin{equation*}
f(\lambda, \xi)=\chi(\xi) \int \frac{\xi^{\frac{1}{2} q} x^{k}}{\xi^{\frac{1}{2} p}\left(m \xi^{\frac{1}{2}}+x\right)^{r}} e^{i\left(\frac{2}{3} m \xi^{\frac{3}{2}}+\frac{1}{3 m^{2}} x^{3}-x \xi\right)} w(\lambda, x) d x \tag{7.48}
\end{equation*}
$$

where supp $\operatorname{sgn}(m)[1, \infty), w \in \lambda^{\frac{1}{6}} L^{2}\left([1, \infty), H_{(-s)}(\mathbb{R})\right)$, and $q+k+2 s \leq r$. We can rewrite this as

$$
\mathscr{F}^{-1}\left(e^{-i \frac{2}{3} m \xi^{\frac{3}{2}}} f\right)(\lambda, x)=a\left(x, D_{x}\right)^{*}\left(e^{i \frac{1}{3 m^{2}}(\cdot)^{3}} w(\lambda, \cdot)\right)
$$

where $\mathcal{F}: L^{2}(\mathbb{R}) \rightarrow L^{2}(\mathbb{R})$ is the Fourier transform and

$$
\begin{equation*}
\left.a(x, \xi)=\left.\chi(2 \operatorname{sgn}(m) x) \chi(\xi) \frac{x^{k} \xi^{\frac{1}{2} q}}{\left(m \xi^{\frac{1}{2}}+x\right)^{r} \xi^{\frac{1}{2} p}} \in\langle m| \xi\right|^{\frac{1}{2}}+|x|\right\rangle^{-r+k} S^{\frac{1}{2}(q-p)}(\mathbb{R} ; \mathbb{R}) \tag{7.49}
\end{equation*}
$$

To see that $f \in \lambda^{\frac{1}{6}} L^{2}(Z)$ if $w \in \lambda^{\frac{1}{6}} L^{2}\left([1, \infty), H_{(-s)}(\mathbb{R})\right)$, it suffices to check that

$$
\left.\left.a \in\langle m| \xi\right|^{\frac{1}{2}}+|x|\right\rangle^{-2 s} S_{1,0}^{0}(\mathbb{R}, \mathbb{R}) \Longrightarrow e^{-i \frac{1}{3 m^{2}} x^{3}} a(x, D): L^{2}(\mathbb{R}) \longrightarrow H_{(s)}(\mathbb{R}), s \geq 0
$$

Since this is clear for $s=0$ it is enough to establish the mapping property for $s \in \mathbb{N}$ as it will then follow by interpolation. Thus let $s=n$ and we need to verify that $D_{x}^{n}\left(\exp \left(-i \frac{1}{3 m^{2}} x^{3}\right) u\right) \in L^{2}(\mathbb{R})$. This however is easy as $|x|^{l} a \in S_{1,0}^{-n+l / 2}$ so that $x^{l} a(x, D): L^{2}(\mathbb{R}) \rightarrow H_{(n-l / 2)}(\mathbb{R})$.

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We will show, using (7.46), that the applications of $D_{\xi}$ and $\xi D_{\xi}$ do not change the form of $f$. To do that we first need to strengthen (7.46) slightly:

$$
\begin{array}{r}
\left(\lambda D_{\lambda}\right)^{k_{0}}\left(x D_{x}\right)^{k_{1}} D_{x}^{k_{2}} w \in \lambda^{\frac{1}{6}} L^{2}\left([1, \infty) ; H_{(-s)}(\mathbb{R})\right)  \tag{7.50}\\
k_{0}+k_{1}+\frac{2}{3} k_{2} \leq 2 l+\frac{2}{3} s, k_{1}+k_{0} \leq 2 l .
\end{array}
$$

The optimal choice of $s$ is in $\frac{1}{2} \mathbb{N}_{0}$ and for $s \in \mathscr{G} \tau_{0}$, (7.50) is clear. For $s \in \mathbb{N}_{0}+\frac{1}{2}$ it is derived from the iteration of the following statement

$$
\begin{align*}
& V^{i} u, D_{x}^{j} u \in \lambda^{\frac{1}{6}} L^{2}(Z), i \leq 2, j \leq 3 \Longrightarrow  \tag{7.51}\\
& D_{x}^{2} V u \in \lambda^{\frac{1}{6}} L^{2}\left([1, \infty) ; H_{\left(-\frac{1}{2}\right)}(\mathbb{R})\right), V=x D_{x}, \lambda D_{\lambda}
\end{align*}
$$

To see this we take the Fourier transform in $x$ and consider $\left\langle\widehat{V^{2} u}, \widehat{D_{x}^{3} u}\right\rangle$. An integration by parts in $\xi$ or $\lambda$ depending on $V$ (in the case of $\lambda$ this means taking an adjoint in of $D_{\lambda}$ ) and an application of the Cauchy-Schwartz inequality yields $\xi^{\frac{2}{3}} \widehat{V u} \in \lambda^{\frac{1}{6}} L^{2}(Z)$ which gives the right hand side of (7.52).
We can now start with

$$
\begin{aligned}
D_{\xi} f(\lambda, \xi)= & \int\left(m \xi^{\frac{1}{2}}-x\right) e^{i\left(\frac{2}{3} m \xi^{\frac{3}{2}}+\frac{1}{3 m^{2}} x^{3}-x \xi\right)} a(x, \xi) w(\lambda, x) d x+ \\
& \int\left(D_{\xi} a\right)(x, \xi) e^{i\left(\frac{2}{3} m \xi^{\frac{3}{2}}+\frac{1}{3 m^{2}} x^{3}-x \xi\right)} w(\lambda, x) d x
\end{aligned}
$$

The second term is of the desired form, as $D_{\xi} a$ is of the form (7.49). We rewrite the first one as

$$
m^{2} \int\left(m \xi^{\frac{1}{2}}+x\right)^{-1} a(x, \xi)\left[\left(-D_{x}\right) e^{i\left(\frac{2}{3} m \xi^{\frac{3}{2}}+\frac{1}{3 m^{2}} x^{3}-x \xi\right)}\right] w(\lambda, x) d x
$$

where we can integrate by parts introducing new $a$ 's of the form (7.49) and $D_{x} w$. By an interpolation argument (see for instance the proof of (7.52) above) it suffices to discuss $D_{\xi}^{2} f, D_{\xi}^{2}$ with weight $\frac{2}{3}$. Repeating the previous computation we see that $D_{\xi}^{2}$ is a sum of terms of the same form as $f$ with $w$ replaced by $D_{x} w$ or $D_{x}^{2} w$ with $r$ in $a$ (see (7.49)) increased by 1 or 2 respectively. If $w \in \lambda^{\frac{1}{6}} L^{2}\left([1, \infty) ; H_{(-s)}\right)$ then (7.50) shows that $D_{x} w \in \lambda^{\frac{1}{6}} L^{2}\left([1, \infty) ; H_{(-s)}\right)$ and $D_{x}^{2} w \in \lambda^{\frac{1}{6}} L^{2}\left([1, \infty) ; H_{(-s-1)}\right)$, where we increased $k_{2}$ by 1 in agreement with the weight of $D_{\xi}^{2}$. Hence, indeed $D_{\xi}^{2} f$ is a sum of terms of the form (7.48) and $D_{\xi}^{2} f \in \lambda^{\frac{1}{6}} L^{2}(Z)$.
We proceed similarly for $\xi D_{\xi}$. Observing that $\xi D_{\xi} a \in S^{0}(\mathbb{R}, \mathbb{R})$ and $\xi e^{i\left(\frac{1}{3 m^{2}} x^{3}-x \xi\right)}=$ $\left(-D_{x}+m^{-2} x^{2}\right) e^{i\left(\frac{1}{3 m^{2}} x^{3}-x \xi\right)}$, we have

$$
\begin{aligned}
& \xi D_{\xi} f(\lambda, \xi)=\int \xi D_{\xi} a(x, \xi) w(\lambda, x) e^{i\left(\frac{2}{3} m \xi^{\frac{3}{2}}+\frac{1}{3 m^{2}} x^{3}-x \xi\right)} d x+ \\
& \quad m^{2} \int\left(m \xi^{\frac{1}{2}}+x\right)^{-1}\left[\left(-D_{x}\right)\left(-D_{x}+m^{-2} x^{2}\right) e^{i\left(\frac{2}{3} m \xi^{\frac{3}{2}}+\frac{1}{3 m^{2}} x^{3}-x \xi\right)} a(x, \xi)\right] w(\lambda, x) d x
\end{aligned}
$$

and again the first term is harmless. Integration by parts in the second term produces terms with new $a$ 's satisfying (7.49) and the following terms involving differentiation of $w$.

$$
\begin{aligned}
& m^{2} \int\left[\left(\frac{2 i}{\left(m \xi^{\frac{1}{2}}+x\right)^{2}} D_{x} w(\lambda, x)+\frac{x}{m^{2}\left(m \xi^{\frac{1}{2}}+x\right)}\left(x D_{x}\right) w(\lambda, x)+\right.\right. \\
& \left.\left.\frac{1}{\left(m \xi^{\frac{1}{2}}+x\right)} D_{x}^{2} w(\lambda, x)\right) a(x, \xi)+\frac{2 D_{x} a(x, \xi)}{m \xi^{\frac{1}{2}}+x} D_{x} w(\lambda, x)\right] e^{i\left(\frac{2}{3} m \xi^{\frac{3}{2}}+\frac{1}{3 m^{2}} x^{3}-x \xi\right)} d x .
\end{aligned}
$$

All the terms except for the third one are of the desired form (since we used the generators in (7.46) with weights $\leq 1$ ). To maintain the order of regularity we use (7.50) again so that if $w \in \lambda^{\frac{1}{6}} L^{2}\left([1, \infty), H_{(-s)}(\mathbb{R})\right)$ then $D_{x}^{2} w \in \lambda^{\frac{1}{6}} L^{2}\left([1, \infty), H_{\left(-s-\frac{1}{2}\right)}\right)$. This decreases $2 l$ by 1 in (7.50) so the order in the filtration is preserved. Since $r$ is increased by $1,(7.48)$ is preserved and $\xi D_{\xi} f \in \lambda^{\frac{1}{6}} L^{2}(Z)$.
To see that $(1-\chi(\xi)) \widehat{v}_{+}(\lambda, \xi) \in \mathcal{S}_{k}^{\alpha} L^{2}(Z)$ we could again analyse the integral (7.47) but it also follows from the uniform ellipticity of $D_{\xi}^{2}-m^{2} \xi$ for $\xi<0$. In fact, if

$$
\operatorname{supp} b \subset(-\infty, 1), \quad\left(\lambda D_{\lambda}\right)^{k_{0}}\left(D_{\xi}^{2}-m^{2} \xi\right)^{k_{1}} b(\lambda, \xi) \in \lambda^{\alpha} L^{2}(Z), \quad k_{0}+\frac{2}{3} k_{1} \leq 2 l
$$

then $b \in \mathscr{S}_{2 l}^{\alpha} L^{2}(Z)$ since

$$
\left\langle\left(D_{\xi}^{2}-m^{2} \xi\right) b, b\right\rangle_{L^{2}\left(\mathbb{R}_{x}\right)} \geq\left\|D_{\xi} b\right\|_{L^{2}\left(\mathbb{R}_{x}\right)}^{2}+m^{2}\left\||\xi|^{\frac{1}{2}} b\right\|_{L^{2}\left(\mathbb{R}_{x}\right)}^{2}-C\|w\|_{L^{2}\left(\mathbb{R}_{x}\right)}^{2}
$$

Note that using an interpolation argument as in the proof of (7.52) we obtain the correct weights for the operators defining $\mathfrak{S}_{2 l}^{\alpha} L^{2}(z)$.
Thus the desired decomposition follows by taking

$$
g(\lambda, \xi)=\widehat{v}_{-}(\lambda, \xi)+(1-\chi(\xi)) \widehat{v}_{+}(\lambda, \xi), \quad f(\lambda, \xi)=e^{\frac{2}{3} i m \xi^{\frac{3}{2}}} \chi(\xi) \widehat{v}_{+}(\lambda, \xi)
$$

We now want to prove the converse, that is, to show that the decomposition on the right hand side of (7.44) implies that $v \in I_{k} L_{\nu_{\frac{1}{6}}}^{2} L^{2}\left(Z, \Xi_{m}\right)$ where as before $\Xi_{m} \stackrel{\text { def }}{=} \Xi_{|m|}^{\mathrm{sgn}(m)}, m \neq 0$. To check the stability under

$$
\lambda D_{\lambda}, D_{x}-m^{-2} x^{2}, x\left(D_{x}-m^{-2} x^{2}\right)
$$

with weights $1, \frac{2}{3}$ and 1 respectively, we move to the Fourier transform side in $x$ and require that

$$
\begin{equation*}
\left(\lambda D_{\lambda}\right)^{k_{0}}\left(D_{\xi}^{2}-m^{2} \xi\right)^{k_{1}}\left[\left(D_{\xi}^{2}-m^{2} \xi\right) D_{\xi}\right]^{k_{2}}\left(\mathcal{F}_{2} \circ W\right) u(\lambda, \xi) \in \lambda^{1 / 6} L^{2}(Z) \tag{7.52}
\end{equation*}
$$

with $k_{0}+\frac{2}{3} k_{1}+k_{2} \leq k=2 l$.
The functions in $\mathscr{S}_{2 l}^{1 / 6} L^{2}(Z)$ satisfy (7.52) and we need to look at $e^{-\frac{2}{3} i m \xi^{\frac{3}{2}}} S_{2 l}^{1 / 6} L^{2}(Z)$. Commutation through the oscillatory factor shows that it suffices to establish the stability of the elements of $S_{2 l}^{1 / 6} L^{2}(Z)$ under $\langle\xi\rangle^{\frac{1}{2}} D_{\xi}$ with the weight $\frac{2}{3}$, i.e.,

$$
f \in S_{2 l}^{1 / 6} L^{2}(Z) \Longrightarrow\left(\langle\xi\rangle^{\frac{1}{2}} D_{\xi}\right)^{p}\left(\lambda D_{\lambda}\right)^{k_{0}}\left(\xi D_{\xi}\right)^{k_{1}} D_{\xi}^{k_{2}} f \in \lambda^{1 / 6} L^{2}(Z)
$$

$\frac{2}{3} p+k_{0}+k_{1}+\frac{1}{3} k_{2} \leq 2 l$.
From this, one sees that it suffices to have the stability under $\langle\xi\rangle^{\frac{1}{2}} D_{\xi}^{2}$ with weight 1 which follows from the stability under $\xi D_{\xi}$ and $D_{\xi}^{3}$ both with weight 1 and the inequality

$$
\begin{aligned}
& \int_{0}^{\infty} t\left|h^{\prime}(t)\right|^{2} d t \leq\left(\int_{0}^{\infty}\left|h^{\prime \prime}(t)\right|^{2} d t\right)^{\frac{1}{2}}\left(\int_{0}^{\infty} t^{2}\left|h^{\prime}(t)\right|^{2} d t\right)^{\frac{1}{2}} \leq \\
& \frac{1}{2}\left(\int_{0}^{\infty}\left|h^{\prime \prime}(t)\right|^{2} d t+\int_{0}^{\infty} t^{2}|h(t)|^{2} d t\right), \quad h \in C_{0}^{\infty}\left(\mathbb{R}_{+}\right)
\end{aligned}
$$

We now need to check the essential support property required in the definition of $I_{k} L_{\nu_{\frac{1}{6}}}^{2}\left(Z, \Xi_{m}\right)$. It is easily satisfied for $\mathcal{F}^{-1} \mathcal{S}_{k}^{1 / 6} L^{2}$ and it remains to consider

$$
v_{+}(x)=\int \chi(2 \xi) e^{-\frac{2}{3} i m \xi^{\frac{3}{2}}+i \xi x} f(\lambda, \xi) d \xi
$$

and to check the stability under $x \cdot \chi(-\operatorname{sgn}(m) x)$ with weight $\frac{1}{3}$. The argument is similar to the analysis of (7.47) and we shall present only the first induction step. Thus we write

$$
x \cdot \chi(-\operatorname{sgn}(m) x) u_{+}(x)=\int \chi(2 \xi) \chi(-\operatorname{sgn}(m) x) \frac{x}{x-m \xi^{\frac{1}{2}}}\left[D_{\xi} e^{i\left(-\frac{2}{3} i m \xi^{\frac{3}{2}}+\xi x\right)}\right] f(\xi) d \xi
$$

where we integrate by parts. If $D_{\xi} f \in \lambda^{1 / 6} L^{2}(Z)$ then so is the left hand side as

$$
a(x, \xi)=\chi(2 \xi) \chi(-\operatorname{sgn}(m) x)\left(x-m \xi^{\frac{1}{2}}\right)^{-1} x \in S^{0}(\mathbb{R}, \mathbb{R})
$$

and we apply $a\left(x, D_{x}\right)$ to $\mathcal{F}^{-1}\left(e^{-\frac{2}{3} i m(\bullet)^{\frac{3}{2}}} D_{\xi} f\right)^{\vee}(x)$. The obvious inductive argument concludes the proof of the lemma.

To prove (7.41) we decompose $S$ into the elliptic and hyperbolic components:

$$
\begin{equation*}
S=S_{e}+S_{h}, S_{e} v(\lambda, x, y)=\frac{1}{(2 \pi)^{\frac{2}{3}}} \int \psi\left(\lambda^{-\frac{2}{3}} \xi\right) \phi(-\xi-y) \frac{A_{+}(-\xi-y)}{A_{+}(-\xi)} \widehat{v}(\xi) e^{i \xi x} d \xi \tag{7.53}
\end{equation*}
$$

where $\phi \in C^{\infty}(\mathbb{R},[0,1])$

$$
\phi(t)= \begin{cases}1 & t>2 \\ 0 & t<1\end{cases}
$$

We start with the elliptic region estimate:
Lemma 7.14. - If $k \in 2 \mathbb{N}_{0}$ and $k_{0}+\frac{2}{3}\left(k_{1}+k_{4}\right)+\frac{1}{3}\left(k_{2}+k_{3}\right) \leq k, \quad m \neq 0$, then $\left(\lambda D_{\lambda}\right)^{k_{0}} D_{x}^{k_{1}} x^{k_{2}} D_{y}^{k_{3}} y^{k_{4}} S_{e}=\mathcal{O}(1): I_{k} L_{\nu_{\alpha}}^{2}\left(Z ; \Xi_{m}^{ \pm}\right) \longrightarrow \lambda^{\alpha} L^{2}\left([1, \infty)_{\lambda} \times \mathbb{R}_{+y} ; H_{\left(\frac{1}{4}\right)}\left(\mathbb{R}_{x}\right)\right)$

Proof. - We again consider $\widetilde{S}_{e}$ defined in the same way as $\widetilde{S}$ in the proof of (7.42) and, using the asymptotic expansions of the Airy functions on $\mathbb{R}_{+}$, we write (omitting d)

$$
\widetilde{S}_{e} \widehat{v}(y, \xi)=\psi\left(\lambda^{-\frac{2}{3}} \xi\right) e^{-\frac{2}{3}\left((-\xi)^{\frac{3}{2}}-(-\xi-y)^{\frac{3}{2}}\right)} w_{1}(-\xi-y) w_{2}(-\xi) \widehat{v}(\xi)
$$

where $w_{i} \in S^{(-1)^{i} / 4}(\mathbb{R})$, supp $w_{i} \subset[1, \infty)$, (compare chapter 9 of [49]). From Lemma 7.13 we conclude that

$$
g(\xi)=\phi(-2 \xi) \widehat{v}(\xi) \in \mathfrak{S}_{k}^{\alpha} L^{2}(Z)
$$

and we want to examine the application of the operators falling on $S_{e}$ in (7.54): $\xi, D_{\xi}, D_{y}, y$ with the weights $\frac{2}{3}, \frac{1}{3}, \frac{1}{3}, \frac{2}{3}$, respectively. From the definition of $\mathfrak{s}_{k}^{\alpha}$, the application of $\lambda D_{\lambda}, \xi$ and $D_{\xi}$ to $g$ is allowed and thus we only need to examine the effect of the last three operators in (7.54) on the phase. Thus

$$
\begin{aligned}
\partial_{\xi}^{k_{2}} \partial_{y}^{k_{3}} y^{k_{4}} \widetilde{S}_{e} \widehat{v}(y, \xi) \equiv & \left((-\xi)^{\frac{1}{2}}-(-\xi-y)^{\frac{1}{2}}\right)^{k_{2}}\left(-(\xi-y)^{\frac{1}{2}}\right)^{k_{3}} y^{k_{4}} \times \\
& e^{-\frac{2}{3}\left((-\xi)^{\frac{3}{2}}-(-\xi-y)^{\frac{3}{2}}\right)} w_{1}(-\xi-y) w_{2}(-\xi) g(\xi) \\
\equiv & \left(\frac{(-\xi-y)^{\frac{1}{2}}}{(-\xi)^{\frac{1}{2}}+(1+y)^{\frac{1}{2}}}\right)^{k_{3}}(1+y)^{k_{4}+\frac{1}{2} k_{3}} \gamma(y, \xi)^{k_{2}} \times \\
& e^{-\frac{2}{3}\left((-\xi)^{\frac{3}{2}}-(-\xi-y)^{\frac{3}{2}}\right)} w_{1}(-\xi-y) w_{2}(-\xi) g(\xi),
\end{aligned}
$$

where $\gamma(y, \xi)=(-\xi)^{\frac{1}{2}}-(-\xi-y)^{\frac{1}{2}}$ and where $\equiv$ means that the expressions differ by terms which can be treated by the analysis for lower $k$. In fact, for $\partial_{\xi}$ and $y$ this is immediate, while for $\partial_{y}$ falling on the phase we write

$$
(-\xi-y)^{\frac{1}{2}}=\frac{(-\xi-y)^{\frac{1}{2}}}{(-\xi)^{\frac{1}{2}}+(1+y)^{\frac{1}{2}}}(1+y)^{\frac{1}{2}}+\frac{(-\xi-y)^{\frac{1}{2}}}{(-\xi)^{\frac{1}{2}}+(1+y)^{\frac{1}{2}}}(-\xi)^{\frac{1}{2}}
$$

with $(-\xi)^{\frac{1}{2}}$ in the second term absorbed into $g$ (with weight $\frac{1}{3}$ ) - it is allowed as $g \in \mathfrak{J}^{\alpha} L^{2}(Z)$. To proceed we need the following

Lemma 7.15. - Let $f \in C^{\infty}\left(\mathbb{R}_{+} \times[1, \infty)\right)$ satisfy
i) $\left|f_{y}^{(k)}(0, t)\right| \leq C_{k}$
ii) $\left|f_{y}^{(k)}(y, t)\right| \leq C_{k} t^{\alpha}$, for some $\alpha \geq 0$
iii) $\operatorname{supp} f \subset\{(y, t): y \geq 0, t \geq 1+y\}$.

Then

$$
\begin{equation*}
\int_{0}^{\infty} e^{-\frac{4}{3}\left(t^{\frac{3}{2}}-(t-y)^{\frac{3}{2}}\right)} f(y, t) d y=\mathcal{O}\left(t^{-\frac{1}{2}}\right) \tag{7.55}
\end{equation*}
$$

Proof. - We can obtain asymptotic expansion for the integral (7.55) using integration by parts:

$$
\begin{aligned}
& \int_{0}^{\infty} e^{-\frac{4}{3}\left(t^{\frac{3}{2}}-(t-y)^{\frac{3}{2}}\right)} f(y, t) d y= \\
& +\frac{1}{2} t^{-\frac{1}{2}} f(0, t)+\int_{0}^{\infty} e^{-4 / 3\left(t^{\frac{3}{2}}-(t-y)^{\frac{3}{2}}\right)}\left[-\frac{1}{4}(t-y)^{-\frac{3}{2}} f(y, t)+\frac{1}{2}(t-y)^{-\frac{1}{2}} f_{y}^{\prime}(y, t)\right] d y
\end{aligned}
$$

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which continued sufficiently many times gives, in view of (i), a term $\mathcal{O}\left(t^{-\frac{1}{2}}\right)$ and a sum of terms of the form

$$
\int_{0}^{\infty} e^{-4 / 3\left(t^{\frac{3}{2}}-(t-y)^{\frac{3}{2}}\right)}(t-y)^{-\frac{k_{1}}{2}} f^{\left(k_{2}\right)}(y, t) d y
$$

with $k_{1}$ large. Using ii) and iii) we can estimate this integral by

$$
\begin{aligned}
& C_{k_{2}} t^{\alpha} e^{-4 / 3 t^{\frac{3}{2}}} \int_{0}^{t-1} e^{4 / 3(t-y)^{\frac{3}{2}}}(t-y)^{-k_{2} / 2} d y \leq C_{k_{2}} t^{\alpha} \int_{0}^{\frac{1}{2} t}(t-y)^{-k_{1} / 2} d y+ \\
& C_{k_{2}} t^{\alpha} e^{-4 / 3\left(1-2^{-\frac{3}{2}}\right) t^{\frac{3}{2}}}=\mathfrak{O}\left(t^{-\frac{1}{2}}\right)
\end{aligned}
$$

if $k_{1} \geq 2 \alpha+3$.
Since $(1+y)^{1 / 2} \leq(-\xi)^{1 / 2}$ for $y \geq 0,-\xi \geq-\xi-y \geq 1$, we can apply Lemma 7.15 with $f(y,-\xi)$ given by

$$
\left[\left(\frac{(-\xi-y)^{\frac{1}{2}}}{(-\xi)^{\frac{1}{2}}+(1+y)^{\frac{1}{2}}}\right)^{k_{3}}(1+y)^{k_{4}+\frac{1}{2} k_{3}} \gamma(y, \xi)^{k_{2}}(-\xi)^{1 / 4} w_{1}(-\xi-y)\right]^{2},
$$

$\alpha=k_{4}+\frac{1}{2}\left(k_{2}+k_{3}\right)$ to obtain from the above discussion

$$
\begin{aligned}
& \iint_{0}^{\infty}\langle\xi\rangle^{\frac{1}{2}}\left|\left(\lambda D_{\lambda}\right)^{k_{0}} D_{x}^{k_{1}} x^{k_{2}} D_{y}^{k_{3}} y^{k_{4}} \widetilde{S}_{e} \widehat{v}(y, \xi)\right|^{2} \lambda^{-2 \alpha} d y d \xi d \lambda \\
& \leq C \sum_{k_{0}^{\prime} \leq k_{0}, k_{1}^{\prime} \leq \frac{1}{2} k_{3}+k_{1}, k_{2}^{\prime} \leq k_{2}} \int\left|\left(\lambda D_{\lambda}\right)^{k_{0}^{\prime}} \xi^{k_{1}^{\prime}} D_{\xi}^{k_{2}^{\prime}} g(\xi)\right|^{2} d \xi \lambda^{-2 \alpha} d \lambda \\
& \leq C \sum_{k_{0}^{\prime}+\frac{2}{3} k_{1}^{\prime}+\frac{1}{3} k_{2}^{\prime} \leq k_{0}+\frac{1}{3} k_{3}+\frac{2}{3} k_{1}+\frac{1}{3} k_{2}} \int\left|\left(\lambda D_{\lambda}\right)^{k_{0}^{\prime}} \xi^{k_{1}^{\prime}} D_{\xi}^{k_{2}^{\prime}} g(\xi)\right|^{2} d \xi \lambda^{-2 \alpha} d \lambda \\
& \leq C\|g\|_{\delta_{k}^{\alpha} L^{2}(Z)} \leq C\|v\|_{I_{k} L_{\nu_{\alpha}}^{2}\left(Z ; \Xi_{m}^{ \pm}\right)}^{2}
\end{aligned}
$$

where the last inequality is a consequence of Lemma 7.13.
We factorize the hyperbolic term defined in (7.53) as follows:

$$
S_{h}=G \circ \tilde{\mathbb{Q}}_{+}^{-1}
$$

where

$$
\begin{align*}
& G v(\lambda, x, y)=\frac{1}{2 \pi} \int \psi\left(\lambda^{-\frac{2}{3}} \xi\right) \Psi_{+}(\tau) \widehat{v}(\xi) e^{-\frac{1}{3} i \tau^{3}+i \tau(\xi+y)+i x \xi} d \tau d \xi  \tag{7.56}\\
& \Psi_{+}(\tau) e^{-\frac{1}{3} i \tau^{3}}=\frac{1}{2 \pi} \int\left((1-\phi) A_{+}\right)(-s) e^{i s t} d s, \quad \Psi_{+} \in S^{0}(\mathbb{R})
\end{align*}
$$

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see [49], Sect. 4. The other term $\widetilde{\mathscr{C}}_{+}^{-1}$ is the non-homogeneous version of the $\mathfrak{Q}_{+}^{-1}$ operator:

$$
\begin{align*}
& \left(\mathfrak{Q}_{+}^{-1} u\right)^{\wedge}\left(\xi_{1}, \xi_{2}\right)=\phi\left(\xi_{1}\right) \psi\left(\xi_{1}^{-1} \xi_{2}\right)\left(A_{+}\left(-\xi_{1}^{-\frac{1}{3}} \xi_{2}\right)\right)^{-1} \widehat{u}\left(\xi_{1}, \xi_{2}\right), u \in \mathcal{G}^{\prime}\left(\mathbb{R}^{2}\right) \\
& \left(\widetilde{\mathscr{a}}_{+}^{-1} v\right)^{\wedge}(\xi)=A_{+}(-\xi)^{-1} \widehat{v}(\xi), \quad v \in \mathfrak{G}^{\prime}(\mathbb{R}) \tag{7.57}
\end{align*}
$$

The more delicate part of the argument involves $\mathbb{Q}_{+}^{-1}$. This operator is associated to a canonical relation with boundary $\mathfrak{e}_{0}^{-}$(see (7.26)) which can be regarded as a singular canonical transformation.

The following proposition, which may be of independent interest, quantizes the Lagrangian mapping properties of $\mathfrak{C}_{0}^{-}$:

$$
\mathfrak{e}_{0}^{-} \Lambda_{m}^{ \pm}=\Lambda_{m \mp 1}^{ \pm}, \quad m \neq 0, \quad \mathfrak{C}_{0}^{-} \Lambda_{0}=\Lambda_{1}^{-}
$$

Proposition 7.16. - Let us define the following marked Lagrangian varieties:

$$
\begin{aligned}
\mathfrak{L}_{p} & =\left\{\Lambda_{|p|}^{\operatorname{sgn}(p)}, 2\right\}, \\
\mathfrak{L}_{0} & =\left\{\Lambda_{0}, J_{\Lambda_{0} \cap \Lambda_{1}}^{1} \Lambda_{1}, 2\right\}
\end{aligned}
$$

where $\Lambda_{m}^{ \pm}, \Lambda_{0}$ are given by (7.28), and $\mathfrak{C}_{0}^{-} \mathfrak{L}_{p}=\mathscr{L}_{p-1}$. Then

$$
\mathfrak{Q}_{+}^{-1}: I_{k} H_{(s)}\left(\mathbb{R}^{2} ; \mathfrak{£}_{p}\right) \longrightarrow I_{k} H_{\left(s-\frac{1}{6}\right)}\left(\mathbb{R}^{2}, \mathfrak{£}_{p-1}\right)
$$

The proof follows from Lemma 7.10 and
Lemma 7.17. - The multiplier $\widetilde{\mathfrak{Q}}_{+}^{-1}$ defined by (7.57) satisfies

$$
\begin{gathered}
\left\langle D_{x}\right\rangle^{-1 / 4} \widetilde{\mathfrak{Q}}_{+}^{-1}: I_{k} L_{\nu_{\alpha}}^{2}\left(Z, \Xi_{m}^{ \pm}\right) \longrightarrow I_{k} L_{\nu_{\alpha}}^{2}\left(Z, \Xi_{m-1}^{ \pm}\right), m>0 \\
\left\langle D_{x}\right\rangle^{-1 / 4} \widetilde{\mathbb{Q}}_{+}^{-1}: I_{k} L_{\nu_{\alpha}}^{2}\left(Z, \Xi_{0}\right) \longrightarrow I_{k} L_{\nu_{\alpha}}^{2}\left(Z, \Xi_{1}^{-}\right)
\end{gathered}
$$

Proof. - This is an easy consequence of Lemma 7.13 and the asymptotic properties of $A_{+}: A_{+}^{-1}=G(\xi)+e^{\frac{2}{3} i \xi^{\frac{3}{2}}} F(\xi), G \in \mathcal{S}(\mathbb{R}), F \in S^{1 / 4}(\mathbb{R}), \operatorname{supp} F \subset[1, \infty)$.

Proof. - Proof of Proposition 7.16 Unless $p=0$ or $p-1=0$ the statement follows immediately from the interpolation between the even indexed neighbours. In the special cases, the proposition is equivalent to

$$
\begin{align*}
& \mathfrak{Q}_{+}^{-1}\left\langle D_{x}\right\rangle^{-s}: I_{k} L^{2}\left(\mathbb{R}^{2}, \mathfrak{Q}_{0}\right) \longrightarrow I_{k} H_{\left(s-\frac{1}{6}\right)}\left(\mathbb{R}^{2}, \mathscr{L}_{-1}\right), p=0  \tag{7.58}\\
& \left\langle D_{x}\right\rangle^{s-\frac{1}{6}} \mathfrak{Q}_{+}^{-1}: I_{k} H_{(s)}\left(\mathbb{R}^{2}, \mathfrak{L}_{1}\right) \longrightarrow I_{k} L^{2}\left(\mathbb{R}^{2}, \mathfrak{Q}_{0}\right), p=1 \tag{7.59}
\end{align*}
$$

In fact, Definition 4.5 immediately implies that for $s \in \mathbb{Z}$,

$$
u \in I_{k} H_{(s)}\left(\mathbb{R}^{2}, \mathfrak{I}_{0}\right) \Leftrightarrow D_{x}^{\alpha} u \in I_{k} L^{2}\left(\mathbb{R}^{2}, \mathfrak{L}_{0}\right),|\alpha| \leq s \Leftrightarrow\left\langle D_{x}\right\rangle^{s} u \in I_{k} L^{2}\left(\mathbb{R}^{2}, \mathfrak{L}_{0}\right)
$$

so that the case of general $s$ follows interpolation. On the other hand Lemma 4.4 shows that $H_{(k)}^{b}\left(Y_{2}\right) \stackrel{\beta_{*}}{\leftrightarrow} I_{k} L^{2}\left(\mathbb{R}^{2}, \mathscr{L}_{0}\right)$, so that by Proposition B. 2 the right hand side is an interpolation space in $k$. Hence (7.58) holds for all $k$.

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Lemma 7.18. - Let $G$ be defined by (7.56). Then for $k$ even

$$
\begin{equation*}
\psi\left(\lambda^{-\frac{1}{3}} x\right) G: I_{k} L_{\nu_{\alpha}}^{2}\left(Z, \Xi_{m}^{ \pm}\right) \longrightarrow I_{k} L_{\nu_{\alpha+\frac{1}{6}}^{+}}^{2}\left(Z_{+}, \widetilde{\Xi}_{m}\right) \tag{7.60}
\end{equation*}
$$

and in particular

$$
\begin{equation*}
\psi\left(\lambda^{-\frac{1}{3}} x\right) G\left\langle D_{x}\right\rangle^{1 / 4}: I_{k} L_{\nu_{\alpha}}^{2}\left(Z, \Xi_{m}^{ \pm}\right) \longrightarrow I_{k} L_{\nu_{\alpha+\frac{1}{3}}^{+}}^{2}\left(Z_{+}, \widetilde{\Xi}_{m}\right) \tag{7.61}
\end{equation*}
$$

Proof. - By taking the Fourier transform in $y$ and applying Young's inequality to eliminate $\psi$ we only need to consider $\lambda^{-1 / 6} \psi\left(\lambda^{-\frac{1}{3}} x\right) \Psi_{+}(\eta) e^{-\frac{1}{3} i \eta^{3}} v(x+\eta)$ and the operators (see definition of $I_{k} L^{2}\left(Z_{+}, \widetilde{\Xi}_{m}\right)$ )

$$
\eta^{2}+D_{\eta}-D_{x}, x+\eta,(x+\eta) D_{x}, \quad m=0
$$

with multiplicities $\frac{2}{3}, \frac{1}{3}, 1$ and

$$
\eta^{2}+D_{\eta}-D_{x}, D_{x}-m^{-2}(x+\eta)^{2},(x+\eta)\left(D_{x}-m^{-2}(x+\eta)^{2}\right), m \neq 0
$$

with multiplicities $\frac{2}{3}, \frac{2}{3}, 1$.
The corresponding assumptions on $v$ gives stability under the operators (see definitions of $I_{k} L^{2}\left(Z, \Xi_{m}^{ \pm}\right): x, x D_{x}, m=0$, with weight $\frac{1}{3}$ and 1), and $D_{x}-m^{-2} x^{2}$, $x\left(D_{x}-m^{-2} x^{2}\right), m \neq 0$, with weight $\frac{2}{3}$ and 1 . It follows that we need to consider the norm of

$$
\lambda^{-1 / 6} \psi\left(\lambda^{-\frac{1}{3}} x\right) a(\eta) w(x+\eta), \quad a \in S^{0}(\mathbb{R}), w \in L^{2}(\mathbb{R})
$$

in $L^{2}\left(\mathbb{R}_{x, \eta}^{2}\right)$ and it is easily seen to be bounded by $\|w\|_{L^{2}(\mathbb{R})}$. The mapping property (7.61) follows from the frequency cut-off $\psi\left(\lambda^{-\frac{2}{3}} \xi\right)$ in $G$.

Combining (7.61) and Lemma 7.17 we obtain that

$$
\begin{gathered}
\psi\left(\lambda^{-\frac{1}{3}} x\right) S_{h}: I_{k} L_{\nu_{\alpha}}^{2}\left(Z, \Xi_{m}^{ \pm}\right) \longrightarrow I_{k} L_{\nu_{\alpha+\frac{1}{3}}^{+}}^{2}\left(Z_{+}, \widetilde{\Xi}_{m \mp 1}\right), m>0 \\
\psi\left(\lambda^{-\frac{1}{3}} x\right) S_{h}: I_{k} L_{\nu_{\alpha}}^{2}\left(Z, \Xi_{0}\right) \longrightarrow I_{k} L_{\nu_{\alpha+\frac{1}{3}}^{+}}^{2}\left(Z_{+}, \widetilde{\Xi}_{1}\right)
\end{gathered}
$$

Thus, using Lemma 7.14 and the definition of $I_{k} L_{\nu_{\alpha+\frac{1}{3}}^{+}}^{2}\left(Z_{+}, \widetilde{\Xi}_{m}\right)$ we obtain (7.41) of Proposition 7.12:

$$
S^{\sharp}=\psi\left(\lambda^{-\frac{1}{3}} x\right) \psi\left(\lambda^{-\frac{2}{3}} y\right)\left(S_{e}+S_{h}\right): I_{k} L_{\nu_{\alpha}}^{2}\left(Z, \Xi_{m}^{ \pm}\right) \longrightarrow I_{k} L_{\nu_{\alpha+\frac{1}{3}}}^{2}\left(Z_{+}, \widetilde{\Xi}_{m}^{ \pm}\right), \quad m \neq 0
$$

Finally, Lemmas 7.10 and 7.11 complete the proof of Proposition 7.9. We should remark that we have not used the full power of Lemma 7.14 which shows precisely the gain of regularity in the elliptic region. A little more will be used in the proof of Proposition 7.21 but even there it would have sufficed to have $\lambda^{\alpha+\frac{1}{6}} L^{2}([1, \infty) \times$ $\left.\mathbb{R}_{+} ; L^{2}(\mathbb{R})\right)$ in the right hand side of (7.54).
7.3. The extension of the equivalence of glancing hypersurfaces to $b$-canonical transformations (see Appendix A and references given there) will now be used to go beyond the model case considered above. Thus we have the following generalization of Proposition 7.9:

Proposition 7.19. - Let the Lagrangians $\Lambda_{m}$ be defined by (7.28) and the relation $\mathcal{F}$ by (7.20). Then the diffractive parametrix $\widetilde{T}$ given by (7.2) has the following mapping properties for $k$ even:

$$
\begin{aligned}
\widetilde{T}: & I_{k} L_{c}^{2}\left(\mathbb{R}^{n-1} ; \Lambda_{0}, J_{\Lambda_{0} \cap \Lambda_{2}}^{1} \Lambda_{2}, 2\right)+I_{k} L_{c}^{2}\left(\mathbb{R}^{n-1} ; \Lambda_{2}^{+}, 2\right) \longrightarrow \\
& I_{k}^{b} L_{\mathrm{loc}}^{2}\left(X, \mathscr{F}\left(\Lambda_{1}\right), \mathscr{F}\left(\Lambda_{1}\right) \cap \mathscr{F}\left(\Lambda_{0}\right)\right), \\
\widetilde{T}: & I_{k} L_{c}^{2}\left(\mathbb{R}^{n-1} ; \Lambda_{1}^{-}, 2\right)+I_{k} L_{c}^{2}\left(\mathbb{R}^{n-1} ; \Lambda_{3}^{+}, 2\right) \longrightarrow I_{k}^{b} L_{\mathrm{loc}}^{2}\left(X, \mathcal{F}\left(\Lambda_{0}\right) \sqcup \mathscr{F}\left(\Lambda_{2}\right)\right), \\
\widetilde{T}: & I_{k} L_{c}^{2}\left(\mathbb{R}^{n-1} ; \Lambda_{1}^{+}, 2\right) \longrightarrow I_{k}^{b} L_{\mathrm{loc}}^{2}\left(X, \mathscr{F}\left(\Lambda_{0}\right) \sqcup \mathcal{F}\left(\Lambda_{\sharp}\right)\right) .
\end{aligned}
$$

where $\Lambda_{\sharp} \subset T^{*} \mathbb{R}^{n-1} \backslash 0$ is any $C^{\infty}$ homogeneous Lagrangian tangent to $\Lambda_{1}$ at $\Lambda_{1} \cap \Lambda_{0}$.
Proof. - Let $\widetilde{\chi}: \widetilde{\Gamma} \rightarrow T^{*} \widetilde{X} \backslash 0$ be the equivalence of glancing hypersurfaces used in the construction of the parametrix. By Proposition 1.1 it can be chosen to extend to a $b$-canonical transformation ${ }^{b} \chi: \Gamma \rightarrow{ }^{b} T^{*} X \backslash 0$. The induced boundary canonical transformation $\chi_{\partial}: \Gamma_{0} \rightarrow T^{*} \partial X \backslash 0$ coincides with $\mathcal{I}$, the canonical transformation of the elliptic Fourier Integral Operator $J$ such that for $v \in \mathfrak{E}^{\prime}\left(\mathbb{R}^{n-1}\right)$

$$
P \widetilde{T} v \equiv 0, \quad \widetilde{T} v \upharpoonright_{X_{-}}=0, \quad \widetilde{T} v \upharpoonright_{\partial X}=J v .
$$

By Theorem 1.5 (and Remark 1.4), ${ }^{b} \chi$ can be chosen so that

$$
\begin{equation*}
\left({ }^{b} \chi\right)^{*} p=c\left(\xi_{2}^{2}-x_{3} \xi_{3}^{2}-\xi_{1} \xi_{3}\right), \quad c \in S_{\mathrm{hg}}^{0}\left(\Gamma_{1}\right) \tag{7.62}
\end{equation*}
$$

where we recall that the pull-back by ${ }^{b} \chi$ makes sense as $p \in S_{\mathrm{hg}}^{2,2}\left({ }^{b} T^{*} X \backslash 0\right)$.
Our first goal is to quantize (7.62) using $b$-Fourier Integral Operators. Thus we choose $F \in I_{b}^{0}\left(\mathbb{R}_{+}^{n} \times X ;\left({ }^{b} \chi^{-1}\right)^{\prime}\right)$ and $G \in I_{b}^{0}\left(X \times \mathbb{R}_{+}^{n} ;\left({ }^{b} \chi\right)^{\prime}\right)$ so that $F G-I=$ $E_{1}, G F-I=E_{2}$, with $E_{1} \in \Psi_{b}^{0}\left(\overline{\mathbb{R}}_{+}^{n}\right), E_{2} \in \Psi_{b}^{0}(X)$ satisfying $W F_{b}\left(E_{1}\right) \cap \Gamma=$ $\emptyset, W F_{b}\left(E_{2}\right) \cap{ }^{b} \chi(\Gamma)=\emptyset\left(W F_{b}(E)\right.$ denotes the essential support of the full symbol of $E$ as a $b$-pseudodifferential operator). By the argument used in the proof of Proposition 4.4 (see also [25], III(4.26) and what follows),

$$
\begin{equation*}
F P G=C\left(P_{0}+R\right)+E, \quad R \in \Psi_{b}^{1,2}\left(\overline{\mathbb{R}}_{+}^{n}\right), E \in \Psi_{b}^{-\infty}\left(\overline{\mathbb{R}}_{+}^{n}\right) \tag{7.63}
\end{equation*}
$$

and where $C \in \Psi_{b}^{0}\left(\overline{\mathbb{R}}_{+}^{n}\right)$ is elliptic in $\Gamma$ and can be chosen (by appropriate modification of $F$ and $G$ ) to have the full symbol vanishing outside a neighbourhood of $\Gamma$. By absorbing the term $D_{x_{2}}^{2} R_{-1}, R_{-1} \in \Psi_{b}^{-1}\left(\overline{\mathbb{R}}_{+}^{n}\right)$, in $R$ into $P_{0}$ (by writing $D_{x_{2}}^{2} R_{-1}=$ $\left.P_{0} R_{-1}+\left(x_{2} D_{x_{3}}^{2}+D_{x_{1}} D_{x_{3}}\right) R_{-1}\right)$ we can assume that $R \in \Psi_{b}^{1,1}\left(\overline{\mathbb{R}}_{+}^{n}\right)$.
The following lemma is based on Proposition 1.2 and allows us to eliminate $R$ :
Lemma 7.20. - If $C \in \Psi_{b}^{0}\left(\overline{\mathbb{R}}_{+}^{n}\right)$ is elliptic in $\Gamma$ and $R \in \Psi_{b}^{1,1}\left(\overline{\mathbb{R}}_{+}^{n}\right)$ then there exist $Q_{1}, Q_{2} \in \Psi_{b}^{0}\left(\overline{\mathbb{R}}_{+}^{n}\right)$, elliptic in $\Gamma$ and $E \in \Psi_{b}^{-\infty, 1}\left(\overline{\mathbb{R}}_{+}^{n}\right)$ such that

$$
C\left(P_{0}+R\right)=Q_{1} P_{0} Q_{2}+E
$$

Proof. - It is in fact more convenient to look for $T_{1}$ and $T_{2}$ in $\Psi_{b}^{0}$, elliptic in $\Gamma$ and such that

$$
T_{1}\left(P_{0}+R\right)=P_{0} T_{2}+E, \quad E \in \Psi_{b}^{-\infty, 1}\left(\overline{\mathbb{R}}_{+}^{n}\right)
$$

which can be rewritten as

$$
\begin{equation*}
\left[P_{0}, T\right] \equiv T R+P_{0} B\left(\bmod \Psi_{b}^{-\infty, 1}\right) \tag{7.64}
\end{equation*}
$$

where $T=T_{1}$ and $B=T_{1}-T_{2}$. Thus we are looking for an elliptic operator $T \in \Psi_{b}^{0}$ and an operator $B \in \Psi_{b}^{-1}$ such that (7.64) holds. Identifying $t^{0}={ }^{b} \sigma_{0}(T) \in$ $S^{0}\left({ }^{b} T^{*} \overline{\mathbb{R}}_{+}^{n}\right)$ and $b^{1}={ }^{b} \sigma_{-1}(B) \in S^{-1}\left({ }^{b} T^{*} \overline{\mathbb{R}}_{+}^{n}\right)$ with their pull-backs under $\jmath$ we see that (7.64) implies

$$
\begin{equation*}
\frac{1}{i} H_{p_{0}} t^{0}=t^{0} r+p_{0} b^{1}, \quad r={ }^{b} \sigma_{1,1}(R) \tag{7.65}
\end{equation*}
$$

On the other hand if (7.65) holds and we choose $T^{0}, B^{1}$ so that ${ }^{b} \sigma_{0}\left(T^{0}\right)=t^{0}$ and ${ }^{b} \sigma_{-1}=b^{1}$ then

$$
\left[P_{0}, T^{0}\right]-T^{0}-P_{0} B^{1} \in \Psi_{b}^{0,1}\left(\overline{\mathbb{R}}_{+}^{n}\right)
$$

Proposition 1.7 shows that, indeed, we can find such $t^{0}$ and $b^{1}$, by taking $t^{0}=\exp a$ and $b^{1}=(\exp a) b$, with $a, b$ given there. We then continue inductively, assuming that

$$
E^{j}=\left[P_{0}, T^{0}+\cdots+T^{j}\right]-\left(T^{0}+\cdots T^{j}\right) R-P_{0}\left(B^{1}+\cdots+B^{j+1}\right) \in \Psi_{b}^{-j, 1}\left(\overline{\mathbb{R}}_{+}^{n}\right)
$$

If ${ }^{b} \sigma_{-j, 1}\left(E^{j}\right)=e^{j}$ we again use Proposition 1.7 to solve

$$
\frac{1}{i} H_{p_{0}}{ }^{j+1}=-e^{j}+r t^{j+1}+p_{0} b^{j+1}
$$

since it is equivalent to
$\frac{1}{i} H_{p_{0}} q=-\left(t^{0}\right)^{-1} e^{j}+p s, \quad t^{j+1}=t^{0} q, b^{j+1}=t^{0}\left(q b^{1}+s\right), q \in S^{-j}(\Gamma), s \in S^{-j-1}(\Gamma)$.
The standard argument adapted to the $\Psi_{b}^{\bullet, 1}$ setting concludes the proof.
By combining Lemma 7.20 with (7.63) we see that there exist $F_{1}, F_{2} \in I_{b}^{0}\left(\overline{\mathbb{R}}_{+}^{n} \times\right.$ $\left.X,\left({ }^{b} \chi^{-1}\right)^{\prime}\right)$ elliptic in $\Gamma$ and properly supported such that

$$
\begin{equation*}
F_{1} P=P_{0} F_{2}+E F_{2}, \quad E \in \Psi_{b}^{-\infty, 1}\left(\overline{\mathbb{R}}_{+}^{n}\right) \tag{7.66}
\end{equation*}
$$

where $E$ is also properly supported. We recall from [25], Proposition III.4.18 that the restriction of a $b$-Fourier Integral Operator, $F$, to the boundary is a Fourier Integral Operator $F_{\partial}$ with the canonical relation given by the restriction of that for $F$. Thus, $\left(F_{2}\right)_{\partial} \in I^{0}\left(\mathbb{R}^{n-1} \times X ;\left(\chi_{\partial}^{-1}\right)^{\prime}\right)$. We conclude that $\left(F_{2}\right)_{\partial} \circ J=A \in \Psi^{0}\left(\mathbb{R}^{n-1}\right)$, where $A$ is properly supported and

$$
\begin{equation*}
P_{0} F_{2} \widetilde{T} v=-E F_{2} \widetilde{T} v+K v, \quad\left(F_{2} \widetilde{T} v\right) \upharpoonright_{\partial x}=A v, \quad K: \mathfrak{G}^{\prime}\left(\mathbb{R}^{n-1}\right) \longrightarrow C^{\infty}\left(\overline{\mathbb{R}}_{+}^{n}\right) \tag{7.67}
\end{equation*}
$$

To continue we need to recall two basic facts. If $X$ and $Y$ are manifolds with boundaries, $\mathfrak{e}$, a $b$-canonical graph then,

$$
\begin{array}{rll}
I_{b}^{0}\left(X \times Y ; \mathcal{E}^{\prime}\right) \ni G & : & \mathscr{\ddots}(Y) \longrightarrow \mathscr{\ddots}(X) \\
\Psi_{b}^{-\infty, k}(X) \ni A & : & \mathscr{\ddots}(X) \longrightarrow C^{\infty}(X) .
\end{array}
$$

The space $9 \tau(X)$ is defined in [25] (see also Sect. 18.3 of [14]). The first property is clear from $\operatorname{III}(4.14)$ in [25], while the second one follows from $\mathfrak{g l}(X) \cap \dot{\mathfrak{d}}(X) \subset$ $\mathscr{B}(X) \cap \dot{Q}(X)=C^{\infty}(X)$ (see [25], Proposition II.8.8 and [14], Theorem 18.3.24) and the mapping properties $\Psi_{b}^{\bullet, k}: \mathfrak{G} \tau \mathfrak{\vartheta} \tau, \Psi_{b}^{-\infty, k}: \mathcal{E}^{\prime} \rightarrow \dot{\mathcal{A}}$.
Combining this with (7.67) we conclude that

$$
F_{2} \widetilde{T}=T_{0} A+E_{0}, \quad A \in \Psi^{0}\left(\mathbb{R}^{n-1}\right), \quad E_{0}: \varepsilon_{G}^{\prime}\left(\mathbb{R}^{n-1}\right) \longrightarrow C^{\infty}\left(\overline{\mathbb{R}}_{+}^{n}\right)
$$

The proof is now concluded through applications of Propositions 4.1,4.2,4.4, Proposition 7.9 and the use of

$$
\chi\left(\mathscr{F}_{0}(\Lambda) \cap \Gamma\right) \subset \mathfrak{F}(\Lambda)
$$

which follows immediately from (7.20).
Proof. - Proof of Theorem 7.6 We proceed by showing that the Poisson map $\mathcal{E}: C_{0}^{\infty}\left(\partial X_{+}\right) \rightarrow C^{\infty}(X):$

$$
\begin{equation*}
P\left(\mathcal{G} f \upharpoonright_{X}\right)=0 \text { in } X, \quad \mathcal{E} f \upharpoonright_{\partial X}=f, \quad \mathcal{G} f \upharpoonright_{X_{-}}=0, \tag{7.68}
\end{equation*}
$$

extends for $s \in 2 \mathbb{N}_{0}$ to a continuous map

$$
\begin{equation*}
\mathfrak{E}: J_{s} L_{c}^{2}(\partial X, H) \longrightarrow\left(\beta_{5}\right)_{*} I_{s} L_{\nu_{5}, \mathrm{loc}}^{2} L^{2}\left(\widetilde{X}_{5}, \varsigma_{5}\right) \tag{7.69}
\end{equation*}
$$

Since both spaces in (7.69) are interpolation spaces, the map is then continuous for all values of $s$, in particular for $s \in \mathbb{N}_{0}$. Once we have (7.69) we apply the 'vanishing in the past' Proposition 5.3 to show that the extension of $\mathcal{E} f \upharpoonright_{X}$ to $\left(\beta_{5}\right)_{*} I_{s} L_{\nu_{5}, \mathrm{loc}}^{2} L^{2}\left(\widetilde{X}_{5}, \mathscr{S}_{5}\right)$ can be modified to lie in $J_{s} L_{\text {comp }}^{2}(\tilde{X}, H)$. In fact, Definition 7.4 shows that for $s \in \mathbb{N}_{0}$

$$
v \in J_{s} L_{\text {comp }}^{2}(\tilde{X}, H) \Rightarrow v \upharpoonright_{K_{\delta}^{\partial}} \in \beta_{*}^{\partial} I_{k} L_{\nu_{\partial}}^{2}\left(\beta^{\partial^{*}} K_{\delta}^{\partial} ; \partial \widetilde{X}_{1} \cap \beta^{\partial^{*}} K_{\delta}^{\partial}\right)
$$

Proposition 5.3 and the spacelike property of $K_{\delta}\left(K_{\delta}^{\partial}=\partial X \cap K_{\delta}\right.$ - see Proposition 5.1) gives the second condition in the definition of $J_{k} L_{\text {loc }}^{2}(\widetilde{X}, H)$ (Definition 3.5; in fact we get a stronger condition: $\left.U \upharpoonright_{K_{1}} \in I_{k} L_{\nu}^{2}\left(\widetilde{X}_{4}, \partial \widetilde{X}_{4}\right) \upharpoonright_{K_{1}}\right)$.
show that the extension of $\mathfrak{G} f \upharpoonright_{X}$ can then be modified to lie in $J_{s} L_{\text {loc }}^{2}(\tilde{X}, H)$ without affecting (7.68).

To establish (7.69) we will reduce it to a microlocal statement. In fact, it suffices to prove (7.69) with the left hand side replaced by the right hand side of (7.6) and the right hand side replaced by the left hand side of (6.2), as we can use Proposition 7.7

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and Theorem 6.1, respectively. In that case we consider the full diffractive parametrix $T$ given by (7.15) and will see that for $k$ even

$$
\begin{align*}
T: & I_{k} L_{c}^{2}\left(\partial X ; \Lambda_{00}, J_{\Lambda_{00} \cap \Lambda_{2 i}}^{1} \Lambda_{2 i}, 2\right)+I_{k} L_{c}^{2}\left(\partial X ; \Lambda_{2 i}^{+}, 2\right) \longrightarrow  \tag{7.70}\\
& I_{k}^{b} L_{\mathrm{loc}}^{2}\left(X, \Lambda_{S}, \Lambda_{S} \cap \Lambda_{F}\right), \\
T: & I_{k} L_{c}^{2}\left(\partial X ; \Lambda_{1 i}^{-}, 2\right)+I_{k} L_{c}^{2}\left(\partial X ; \Lambda_{3 i}^{+}, 2\right) \longrightarrow I_{k}^{b} L_{\mathrm{loc}}^{2}\left(X, \Lambda_{F} \sqcup \Lambda_{R}\right),  \tag{7.71}\\
T: & I_{k} L_{c}^{2}\left(\partial X ; \Lambda_{11}^{+}, 2\right) \longrightarrow I_{k}^{b} L_{\mathrm{loc}}^{2}\left(X, \Lambda_{F} \sqcup \Lambda_{R}\right),  \tag{7.72}\\
T: & I_{k} L_{c}^{2}\left(\partial X ; \Lambda_{13}^{+}, 2\right) \longrightarrow I_{k}^{b} L_{\mathrm{loc}}^{2}\left(X, \Lambda_{H} \sqcup \Lambda_{S}\right) . \tag{7.73}
\end{align*}
$$

Recalling (7), we see that, after a suitable microlocalization, (7.70),(7.71),(7.72) and (7.73) are a consequence of Proposition 7.19 and the following mapping properties of $L$ in (7.15):

$$
\begin{array}{rll}
L: & I_{k} L_{c}^{2}\left(\partial X ; \Lambda_{00}, J_{\Lambda_{00} \cap \Lambda_{2 i}}^{1} \Lambda_{2 i}, 2\right)+I_{k} L_{c}^{2}\left(\partial X ; \Lambda_{2 i}^{+}, 2\right) \longrightarrow \\
& & I_{k} L_{c}^{2}\left(\mathbb{R}^{n-1} ; \Lambda_{0}, J_{\Lambda_{0} \cap \Lambda_{2}}^{1} \Lambda_{2}, 2\right)+I_{k} L_{c}^{2}\left(\mathbb{R}^{n-1} ; \Lambda_{2}^{+}, 2\right), \\
L: & I_{k} L_{c}^{2}\left(\partial X ; \Lambda_{1 i}^{-}, 2\right)+I_{k} L_{c}^{2}\left(\partial X ; \Lambda_{3 i}^{+}, 2\right) \longrightarrow \\
& I_{k} L_{c}^{2}\left(\mathbb{R}^{n-1} ; \Lambda_{1}^{-}, 2\right)+I_{k} L_{c}^{2}\left(\mathbb{R}^{n-1} ; \Lambda_{3}^{+}, 2\right), \\
L: & I_{k} L_{c}^{2}\left(\partial X ; \Lambda_{11}^{+}, 2\right) \longrightarrow I_{k} L_{c}^{2}\left(\mathbb{R}^{n-1} ; \Lambda_{1}^{+}, 2\right) \\
L: & I_{k} L_{c}^{2}\left(\partial X ; \Lambda_{13}^{+}, 2\right) \longrightarrow I_{k} L_{c}^{2}\left(\mathbb{R}^{n-1} ; \Lambda_{1}^{+}, 2\right) . \tag{7.77}
\end{array}
$$

The operator $L$ is chosen differently in each case depending on $\chi$ and $\Gamma$ used in the parametrix construction - see the discussion following (7.19). Since $L$ is an elliptic Fourier Integral Operator associated with $\chi_{\partial}$ (see (7.19)) the mapping properties above follow from Propositions 4.1 and 4.2. We then obtain (7.70), (7.71),(7.72) and (7.73) from (7.74),(7.75),(7.76) and (7.77), respectively from Proposition 7.6, once we check, for (7.72), that

$$
\Lambda_{R} \cap \tilde{\Gamma} \supset \mathcal{F}\left(\Lambda_{\sharp} \cap \Gamma\right),
$$

and for (7.73), that

$$
\Lambda_{S} \cap \tilde{\Gamma} \supset \mathcal{F}\left(\Lambda_{\sharp} \cap \Gamma\right)
$$

for some $\Lambda_{\sharp}$ tangent to $\Lambda_{1}$ at $\Lambda_{1} \cap \Lambda_{0}$ (again, with different choices of $\chi$ and $\widetilde{\Gamma}$ in the two cases (7.72) and (7.73)). For (7.72) we can simply take $\Lambda_{\sharp}=\Lambda_{2}=\bigcup_{ \pm} \delta_{0}^{ \pm} \Lambda_{0}$ as then $\mathcal{F}\left(\Lambda_{2} \cap \Gamma\right) \subset \Lambda_{R} \cap \widetilde{\Gamma}$. To construct $\Lambda_{\sharp}$ for (7.73) we observe that if $\mathcal{F}$ and $\mathcal{F}_{1}$ are defined using canonical transformations $\widetilde{\chi}$ and $\widetilde{\chi}_{1}-$ see $(7.20)-$ then $\mathcal{F}_{1}=\mathscr{F} \circ \chi_{\partial}^{-1}\left(\chi_{1}\right)_{\partial}$. If $\left(\chi_{1}\right)_{\partial}$ is chosen so that $\left(\chi_{1}\right)_{\partial}\left(\Lambda_{0} \cap \Gamma\right) \subset \Lambda_{00}$, as it may be by (7.19), then by (7)

$$
\mathcal{F}_{1}\left(\Lambda_{1} \cap \Gamma\right) \subset \Lambda_{S} \cap \widetilde{\chi}(\widetilde{\Gamma})
$$

and we can take $\Lambda_{\sharp}=\chi_{\partial}^{-1}\left(\chi_{1}\right)_{\partial}\left(\Lambda_{1}\right)$. Since $\Lambda_{1}$ is tangent to $\Lambda_{2}$, it remains to check that $\Lambda_{\sharp}$ is tangent to $\Lambda_{2}$, and that follows from the tangency of $\mathcal{F}\left(\Lambda_{\sharp}\right)$ to $\mathscr{F}\left(\Lambda_{2}\right)$, as can be verified using $\mathcal{F}_{0}$ or even more easily $\mathscr{H}$ given by (7.27). By the last part of (7), $\mathscr{F}\left(\Lambda_{2} \cap \Gamma\right) \subset \Lambda_{H_{+}}$and that tangency follows from the third order tangency of the cone and the reflected front $H_{+}$(see chapter 2). This completes the proof of Theorem 7.6 .
7.4. Because of the restriction to the boundary (see Theorem 7.5) we lose $1 / 2$ of the order of regularity when applying Theorem 7.6. To avoid this loss in the final application to conormal regularity of diffracted waves we need

Proposition 7.21. - If $\left(F \cap \bar{X}_{-}\right) \cap \partial X=\emptyset$ and $u \in L^{2}(X)$ satisfies

$$
\begin{equation*}
P u=0 \quad \text { in } X, u \upharpoonright_{\partial X}=0, u \upharpoonright_{X_{-}}=u_{0}, \quad u_{0} \in I_{k} L^{2}\left(X_{-} ; \Lambda_{F}\right) \tag{7.78}
\end{equation*}
$$

then

$$
u \in{ }^{b} I_{k} L_{\mathrm{loc}}^{2}\left(X ; \Lambda_{F} \sqcup \Lambda_{R}\right) .
$$

Proof. - Following the proof of Theorem 7.6 and with the notation of Proposition 7.9, we only need to prove the statement in the model case: $P=P_{0}, \Lambda_{F}=\mathscr{F}_{0}\left(\Lambda_{0}\right)$, $\Lambda_{R}=\mathscr{F}_{0}\left(\Lambda_{2}\right), X=\mathbb{R}_{+}^{n}, X_{-}=\left\{\phi(x)<-\delta, x_{2}>0\right\}$, where for small $|x|$ we can take $\phi(x)=x_{1}+x_{3}$. The assumption that $\bar{X}_{-} \cap F$ is away from the boundary is replaced by its microlocal version:

$$
\begin{equation*}
W F\left(u_{0}\right) \subset \mathfrak{F}_{0}\left(\Gamma_{0}\right) \cap \pi^{-1}\left(X_{-}\right) \tag{7.79}
\end{equation*}
$$

where $\Gamma_{0} \subset T^{*} \mathbb{R}^{n-2}$ is small conic neighbourhood of $(0 ;(1,0, \cdots, 0)$ ).
The lemma below provides the needed characterization of the restriction of the free solution which is already essentially contained in Corollary 5.10 of [35] (see also Sect. 25.3 of [14] and Sect. 3 of [49]).

Lemma 7.22. - If $u_{0}$ satisfies (7.79) and if

$$
P_{0} w=0 \text { in } \mathbb{R}^{3}, w \upharpoonright_{\left(\mathbb{R}_{+}^{3}\right)_{-}}=u_{0}, u_{0} \in I_{k} L_{\mathrm{loc}}^{2}\left(\left(\mathbb{R}^{3}\right)_{-}, \mathscr{F}_{0}\left(\Lambda_{0}\right)\right)
$$

then $w \in I_{k} L_{\mathrm{loc}}^{2}\left(\mathbb{R}^{3} ; \mathfrak{F}_{0}\left(\Lambda_{0}\right)\right)$ and

$$
w \upharpoonright_{\partial \mathbb{R}_{+}^{3}}=\mathfrak{C l i} w_{0}+E u_{0}, w_{0} \in I_{k} H_{(-1 / 6)}\left(\mathbb{R}^{2}, \Lambda_{0}\right),
$$

where $\mathfrak{C l} i$ is the multiplier defined by $A i$ as in (7.57) and $E: L_{\mathrm{loc}}^{2}\left(\mathbb{R}^{3}\right) \rightarrow C^{\infty}\left(\mathbb{R}^{3}\right)$.
Proof. - The first property of $w$ is immediate from the propagation properties of Lagrangian distributions. Condition (7.79) implies that $W F\left(u_{0}\right)$ is contained in a conic neighbourhood of the bicharacteristic for $p_{0}$ through $(0 ;(0,0,1))$ :

$$
x_{2}=x_{1}^{2}, x_{3}=\frac{2}{3} x_{2}^{3}, \xi_{1}=0, \xi_{2}=-x_{1}, \xi_{3}=1
$$

Thus by smoothly cutting off of the initial data (which produces $C^{\infty}$ errors) we can assume that $u_{0} \in C^{\infty}\left(\mathbb{R}_{x_{2}} ; \mathfrak{E}^{\prime}\left(\mathbb{R}_{x_{1}, x_{3}}^{2}\right)\right)$ and consequently that the same is true for $\psi w$, where $\psi \in C^{\infty}\left(\mathbb{R}_{x_{2}}, C_{0}^{\infty}\left(\mathbb{R}_{x_{1}, x_{2}}\right)\right), \psi \equiv 1$ near $\pi\left(\mathscr{F}_{0}\left(\Gamma_{0}\right)\right)$. We also note that $(1-\psi) w \in C^{\infty}\left(\mathbb{R}^{3}\right)$. Taking the Fourier transform $\mathcal{F}$ in $x_{1}$ and $x_{3}$ we obtain

$$
\left(D_{x_{2}}^{2}-\xi_{3}\left(\xi_{3} x_{2}+\xi_{1}\right)\right) \widehat{\psi w}\left(\xi_{1}, x_{2}, \xi_{3}\right)=\mathscr{F}\left(-\left[P_{0}, \psi\right] w\right)\left(\xi_{1}, x_{2}, \xi_{3}\right) \in C^{\infty}\left(\mathbb{R}_{x_{2}} ; \mathscr{\delta}\left(\mathbb{R}_{\xi_{1}, \xi_{3}}^{2}\right)\right)
$$

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and since $\widehat{\psi w}$ is tempered in $\left(\xi_{1}, \xi_{3}\right)$ we conclude that

$$
\phi\left(\xi_{3}\right) \widehat{\psi w}\left(\xi_{1}, x_{2}, \xi_{3}\right)=\phi\left(\xi_{3}\right) A i\left(-\xi_{3}^{-1 / 3}\left(\xi_{1}+\xi_{3} x_{2}\right)\right) \widehat{w}_{0}\left(\xi_{1}, \xi_{3}\right)
$$

where we can neglect the term $\phi$ in the left hand side (see $(7.57))$ as $\mathcal{F}^{-1}((1-\phi) \widehat{\psi w}) \in$ $C^{\infty}$. If $\Gamma_{0}$ is small enough, then $x_{2}>\delta_{1}>2 \max _{\xi \in \Gamma_{0}}\left|\xi_{1} / \xi_{3}\right|$ in $\left(\mathbb{R}^{3}\right)_{-} \cap \operatorname{sing} \operatorname{supp} u_{0}$. Thus, for $(x, \xi) \in \mathcal{F}_{0}\left(\Gamma_{0}\right) \cap \pi^{-1}\left(X_{-} \cap \operatorname{sing} \operatorname{supp} u_{0}\right)$

$$
\begin{aligned}
& \phi\left(\xi_{3}\right) A i\left(-\xi_{3}^{-1 / 3}\left(\xi_{1}+\xi_{3} x_{2}\right)\right)=\xi_{3}^{-1 / 6}\left(x_{2}+\xi_{1} / \xi_{3}\right)^{-1 / 4} \phi\left(\xi_{3}\right) \times \\
& {\left[e^{i \frac{2}{3} \xi_{3}\left(x_{2}+\xi_{1} / \xi_{3}\right)^{3 / 2}} a_{+}\left(\xi_{3}^{2 / 3}\left(x_{2}+\xi_{1} \xi_{3}\right)\right)+e^{-i \frac{2}{3} \xi_{3}\left(x_{2}+\xi_{1} / \xi_{3}\right)^{3 / 2}} a_{-}\left(\xi_{3}^{2 / 3}\left(x_{2}+\xi_{1} \xi_{3}\right)\right)\right]}
\end{aligned}
$$

$a_{ \pm} \in S_{\mathrm{phg}}^{0}(\mathbb{R}), a_{ \pm}(1) \neq 0$. Since the phases are smooth then, we conclude that

$$
\begin{gathered}
W F\left(\mathscr{F}^{-1}\left(\exp \left( \pm \frac{2}{3} i \xi_{3}\left(x_{2}+\xi_{1} / \xi_{3}\right)^{3 / 2}\right) a_{ \pm} \widehat{w}_{0}\right) \upharpoonright_{\left\{x_{2}>\delta_{1},|x|<\delta_{2}\right\}}\right) \subset \\
\mathcal{F}_{0}\left(W F\left(w_{0}\right)\right) \cap\left\{x_{2}>\delta_{1}, \pm x_{1}<0\right\}
\end{gathered}
$$

Hence, up to an error in $C^{\infty}\left(\mathbb{R}_{x_{2}} ; \mathcal{S}\left(\mathbb{R}_{\xi_{1}, \xi_{3}}^{2}\right)\right)$
$\widehat{u}_{0}\left(\xi_{1}, x_{2}, \xi_{3}\right)=\xi_{3}^{-\frac{1}{6}}\left(x_{2}+\xi_{1} / \xi_{3}\right)^{-1 / 4} \phi\left(\xi_{3}\right) e^{i \frac{2}{3} \xi_{3}\left(x_{2}+\xi_{1} / \xi_{3}\right)^{3 / 2}} a_{+}\left(\xi_{3}^{2 / 3}\left(x_{2}+\xi_{1} \xi_{3}\right)\right) \widehat{w}_{0}\left(\xi_{1}, \xi_{3}\right)$.
If $u_{0} \in L_{\text {loc }}^{2}\left(\mathbb{R}_{x_{2}} ; L^{2}\left(\mathbb{R}_{x_{1}, x_{2}}^{2}\right)\right)$ we conclude that $\widehat{w}_{0} \in\left\langle\xi_{3}\right\rangle^{1 / 6} L^{2}\left(\mathbb{R}^{2}\right)$, that is, $w_{0} \in$ $H_{-\frac{1}{6}}$. The desired conclusion is equivalent to the stability condition

$$
\left(\xi_{3} D_{\xi_{1}}\right)^{k_{1}}\left(\xi_{3} D_{\xi_{3}}\right)^{k_{2}} \widehat{w}_{0} \in\langle\xi\rangle^{1 / 6} L^{2}\left(\mathbb{R}^{2}\right), \quad k_{1}+k_{2} \leq k
$$

while we know that

$$
\left(-\xi_{3} D_{\xi_{3}}+2 x_{2} D_{x_{2}}-\xi_{1} D_{\xi_{1}}\right)^{k_{0}}\left(D_{x_{2}}-\xi_{3} D_{\xi_{1}}\right)^{k_{1}} \widehat{u}_{0} \in L_{\mathrm{loc}}^{2}\left(\mathbb{R}_{x_{2}} ; L^{2}\left(\mathbb{R}_{\xi}^{2}\right)\right)
$$

However,

$$
\begin{aligned}
& \left(D_{x_{2}}-\xi_{3} D_{\xi_{1}}\right) \widehat{u}_{0}= \\
& \xi_{3}^{-1 / 6}\left(x_{2}+\xi_{1} / \xi_{3}\right)^{-1 / 4} \phi\left(\xi_{3}\right) e^{i \frac{2}{3} \xi_{3}\left(x_{2}+\xi_{1} / \xi_{3}\right)^{3 / 2}} a_{+}\left(\xi_{3}^{2 / 3}\left(x_{2}+\xi_{1} \xi_{3}\right)\right)\left(-\xi_{3} D_{\xi_{1}}\right) \widehat{w}_{0} \\
& \left(-\xi_{3} D_{\xi_{3}}+2 x_{2} D_{x_{2}}-\xi_{1} D_{\xi_{1}}\right) \widehat{u}_{0} \equiv \xi_{3}^{-1 / 6}\left(x_{2}+\xi_{1} / \xi_{3}\right)^{-1 / 4} \phi\left(\xi_{3}\right) \times \\
& e^{i \frac{2}{3} \xi_{3}\left(x_{2}+\xi_{1} / \xi_{3}\right)^{3 / 2}} a_{+}\left(\xi_{3}^{2 / 3}\left(x_{2}+\xi_{1} \xi_{3}\right)\right)\left(-3 \xi_{3} D_{\xi_{3}}-\xi_{1} D_{\xi_{1}}\right) \widehat{w}_{0}
\end{aligned}
$$

which together with the wave-front assumption (in particular, $\left|\xi_{1}\right| \leq \xi_{3}$ ) concludes the proof.

With the notation of Proposition 7.12, we want to show that

$$
S^{\sharp}\left(W\left(\mathfrak{C i} w_{0}\right)\right) \in I_{k} L_{\nu_{\frac{1}{2}}^{+}}^{2}\left(Z_{+}, \widetilde{\Xi}_{0}\right)+I_{k} L_{\nu_{\frac{1}{2}}^{+}}^{2}\left(Z_{+}, \widetilde{\Xi}_{2}\right),
$$

as then Lemma 7.11 will provide the desired conclusion. Again, we consider the elliptic and hyperbolic components separately and start with the former:

$$
S_{e}\left(W\left(\mathbb{Q} i w_{0}\right)\right)=S_{e}\left(\phi\left(-D_{x}\right) \widetilde{\mathfrak{d}} i W\left(w_{0}\right)\right), \quad \widetilde{\mathfrak{d}} i \circ W \equiv W \circ \mathfrak{Q} i
$$

Since $A i(-\xi)=\mathfrak{O}\left(\langle\xi\rangle^{-N}\right)$ for all $N$ if $\xi \leq 0$, we easily see that $\phi(-\xi) A i(-\xi) \widehat{W w_{0}}(\lambda, \xi) \in$ $\mathscr{S}_{k}^{1 / 3} L^{2}(Z)$. Thus we can use the proof of Lemma 7.14 to conclude that

$$
S_{e}\left(\tilde{\mathscr{q}}_{i} W\left(w_{0}\right)\right) \in I_{k} L_{\nu_{\frac{1}{2}}^{+}}^{2}\left(Z_{+}, \widetilde{\Xi}_{2}\right)
$$

In fact, a much stronger conclusion holds:

$$
\left(\lambda D_{\lambda}\right)^{k_{0}} D_{x}^{k_{1}} x^{k_{2}} D_{y}^{k_{3}} y^{k_{4}} S_{e}\left(\widetilde{\mathbb{q}} i W\left(w_{0}\right)\right) \in \lambda^{1 / 3} L^{2}\left([1, \infty) \times \mathbb{R} ; H_{(1 / 4)}(\mathbb{R})\right) \subset \lambda^{1 / 2} L^{2}\left(Z_{+}\right)
$$

Turning to the hyperbolic component we observe that by writing $A i=-\omega A_{+}-\bar{\omega} A_{-}$, $\omega=\exp (2 \pi i / 3)$, we obtain

$$
S_{h}\left(\widetilde{\mathfrak{q}} i\left(W w_{0}\right)\right)=-\omega G W w_{0}-\bar{\omega} G\left(\widetilde{\left(\frac{\mathfrak{Q}_{-}}{\mathfrak{Q}_{+}}\right)} W w_{0}\right)
$$

Straightforward analogues of Lemmas 7.17 and 7.18 conclude the proof.
We should remark that a minimal amount of additional care would remove the assumption that $k$ is even, but that is irrelevant to us as indicated in the following

Theorem 7.23. - If the assumptions of Proposition 7.21 are satisfied with $k \in \mathbb{N}_{0}$, then there exists a continuous map

$$
E: I_{k} L^{2}\left(X_{;} \Lambda_{F}\right) \longrightarrow J_{k} L_{c}^{2}(\widetilde{X}, H)
$$

such that $E\left(u_{0}\right) \upharpoonright_{X}=u$.
Proof. - We only need to combine Proposition 7.21 with Proposition 5.3, Theorem 6.1 and an interpolation argument as in the proof of Theorem 7.6.

## 8. PROOF OF THE MAIN THEOREM

To give the proof of the propagation theorem for the pseudo-conormal space given by Definition 3.6, two preparatory facts are still needed.

Lemma 8.1. - If $g \in J_{s} L_{\text {loc }}^{2}(\tilde{X}, H)$ then
$W F^{(s)}(g) \subset N^{*} \widetilde{F} \cup N^{*} \widetilde{R} \cup N^{*} \widetilde{S}_{+} \cup T_{\Gamma}^{*} \widetilde{X} \backslash 0 \cup N^{*} L \cup N^{*} B_{+} \cup N^{*} D_{+} \cup N^{*} H$.

Proof. - The statement follows easily from Definition 3.5 if $s \in \mathbb{N}_{0}$. We need to see that it is preserved when the space $J_{k} L_{\text {loc }}^{2}(\tilde{X}, H)$ is interpolated. For convenience we shall denote the union on the right hand side of (8.1) by $\widetilde{\mathscr{L}}$.
Let $\Omega \subset T^{*} \tilde{X} \backslash 0$ be an open conic set such that $\bar{\Omega}$ is disjoint from $\widetilde{\mathscr{L}}$. The condition $W F^{(s)}(g) \subset \widetilde{\mathscr{L}}$ is equivalent to saying that for every such $\Omega$, if $a \in S_{\mathrm{hg}}^{0}\left(T^{*} X \backslash 0\right)$ is elliptic in $\Omega$ and supp $a \cap \widetilde{\varrho}=\emptyset$, then $a(x, D)\langle D\rangle^{s} g \in L_{\text {loc }}^{2}(X)$. Fixing $a$ with that property and defining

$$
W_{s}(\tilde{X}, \Omega)=\left\{u \in L_{\mathrm{loc}}^{2}(X): a(x, D)\langle D\rangle^{s} g \in L_{\mathrm{loc}}^{2}(\tilde{X})\right\}
$$

we see that $W_{s}$ is an interpolation space in $s$ and that $W F^{(s)}(g) \subset \widetilde{\mathfrak{~}}$ if and only if for all $\Omega$ above $g \in W_{s}(\widetilde{X}, \Omega)$. Fixing $\Omega$ we now see that

and the complex interpolation finishes the proof.
To guarantee the singular support condition in Definition 3.6 of $J_{s} L_{\text {loc }}^{2}(X)$ we need the following crude propagation property:

Lemma 8.2. - Let u satisfy

$$
\begin{equation*}
P u=f \text { in } X, u \upharpoonright_{\partial X}=0, u \upharpoonright_{X_{-}}=0 \tag{8.2}
\end{equation*}
$$

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with $f=\tilde{f} \upharpoonright_{X}$, where $W F^{(s)}(\tilde{f})$ is contained in the right hand side of (8.1) If

$$
\text { sing supp }{ }^{(s)}(f) \cap\left(\widetilde{R} \backslash R \cup \widetilde{F} \backslash F \cup \widetilde{S}_{+} \backslash S_{+}\right)=\emptyset
$$

then

$$
\text { sing supp }{ }^{\left(s+\frac{1}{2}\right)}(u) \cap\left(\widetilde{R} \backslash R \cup \widetilde{F} \backslash F \cup \widetilde{S}_{+} \backslash S_{+}\right)=\emptyset
$$

Proof. - We construct the solution $u$ by starting with the interior problem:

$$
\widetilde{P} \tilde{u}=\tilde{f} \text { in } \tilde{X}, \quad \tilde{u} \upharpoonright_{\tilde{X}_{-}}=0
$$

where we choose any strictly hyperbolic extension of $P \upharpoonright_{X}$ to $\tilde{X} \backslash X$ (a fixed extension $P$ was introduced in chapter 2). If we denote the right-hand side of (8.1) by $\widetilde{\mathscr{\varrho}}$ and introduce

$$
\mathfrak{L}=\widetilde{\mathfrak{L}} \backslash \pi^{-1}\left[\left(\widetilde{R} \backslash R \cup \widetilde{F} \backslash F \cup \widetilde{S}_{+} \backslash S_{+}\right) \cap X\right],
$$

then the assumption on $f$ shows that $W F^{(s)}(\tilde{f}) \subset \mathfrak{L}$. If $\widetilde{\mathfrak{L}}_{1}$ is the $\widetilde{P}$-flow-out of $\widetilde{\varrho} \cap\left\{\sigma_{2}(\widetilde{P})=0\right\}$, then the standard propagation result (see [14], Sect. 26.1) gives

$$
W F^{(s+1)}(\tilde{u}) \subset \mathfrak{L} \cup \widetilde{\mathfrak{\Omega}}_{1}
$$

By choosing $\widetilde{P}$ in $\tilde{X} \backslash X$ appropriately, it can be arranged (for instance, by decreasing the speed of propagation in $\tilde{X} \backslash X$ ) that

$$
\begin{equation*}
\left(\widetilde{\mathscr{L}}_{1} \backslash \mathfrak{L}\right) \cap(\widetilde{\mathscr{L}} \backslash \mathfrak{L})=\emptyset, \tag{8.3}
\end{equation*}
$$

and this will be crucial later. At the boundary, all the terms in $\widetilde{\varrho}$ and $\widetilde{\mathscr{I}}_{1}$ are disjoint except over $\Gamma$ where all of $T_{\Gamma}^{*} \widetilde{X} \backslash 0$ is included. Thus

$$
W F^{\left(s+\frac{1}{2}\right)}\left(-\tilde{u} \upharpoonright_{\partial X}\right) \subset \widetilde{\mathfrak{L}}_{1}^{\partial}
$$

where $\widetilde{\mathfrak{L}}_{1}^{\partial}$ comes from the projection of each term in $\widetilde{\mathfrak{L}}_{1}$ to $T^{*} \partial X \backslash 0$. Hence when we solve

$$
P v=0 \text { in } X, v \upharpoonright_{\partial X}=-\tilde{u} \upharpoonright_{\partial X}, v \upharpoonright_{X_{-}}=0,
$$

the propagation of singularities theorem for the diffractive Dirichlet problem (see [37] and [14], Sect. 24.4) shows that

$$
v \in \mathfrak{g l}(X), \quad W F_{b}^{\left(s+\frac{1}{2}\right)}(v) \subset \jmath\left(\widetilde{\mathscr{®}}_{1} \upharpoonright_{T^{*} X \backslash 0}\right)
$$

where $\jmath: T^{*} X \backslash X \rightarrow{ }^{b} T^{*} X \backslash 0$. Putting $u=\tilde{u} \upharpoonright_{X}+v$ we obtain the solution to (8.2) and it is independent of the extension $\widetilde{P}$ chosen. Thus we can take either the fixed extension $P$ (as in chapter 2 and consequently in the definition the extended fronts and $\widetilde{\mathscr{S}}$ ) or $\widetilde{P}$ such that (8.3) holds. This shows that

$$
\operatorname{sing} \operatorname{supp}^{\left(s+\frac{1}{2}\right)}(u) \subset \pi(\widetilde{\mathfrak{L}}) \upharpoonright_{X} \cap \pi\left(\widetilde{\mathscr{L}}_{1}\right) \upharpoonright_{X} \subset \pi(\mathfrak{L}) \upharpoonright_{X}
$$

concluding the proof.

We can finally give the long promised
Theorem 8.3. - If $f \upharpoonright_{X_{-}}=0$ and $u \in L_{l o c}^{2}(X)$ satisfies

$$
P u=f \quad \text { in } X, \quad u \upharpoonright_{\partial X}=0, \quad u \upharpoonright_{X_{-}}=0,
$$

then

$$
f \in J_{s} L_{\mathrm{loc}}^{2}(X) \Longrightarrow u \in J_{s+\frac{1}{2}} L_{\mathrm{loc}}^{2}(X)
$$

Proof. - Definition 3.6 guarantees that for any $H \in \mathcal{R}$ there exists $\tilde{f} \in$ $J_{s} L_{\text {loc }}^{2}(\tilde{X}, H)$ such that $f=\tilde{f} \upharpoonright_{X}$. If we apply Theorem 5.6 to the problem

$$
P u_{1}=\tilde{f} \text { in } \tilde{X}, \quad u_{1} \upharpoonright_{\tilde{X}_{-}}=0
$$

we conclude that $u_{1} \in J_{s}^{1} L_{\text {loc }}^{2}(\tilde{X}, H)$. Theorem 7.5 then implies that

$$
-u_{1} \upharpoonright_{\partial X} \in J_{s+\frac{1}{2}}(\partial X, H)
$$

and we solve the Dirichlet problem with that boundary data:

$$
P u_{2}=0 \text { in } X, \quad u_{2} \upharpoonright_{X_{-}}=0, \quad u_{2} \upharpoonright_{\partial X}=-u_{1} \upharpoonright_{\partial X}
$$

Theorem 7.6 shows that $u_{2} \in J_{s+\frac{1}{2}} L_{\text {loc }}^{2}(\tilde{X}, H) \upharpoonright_{X}$ and $u=\tilde{u} \upharpoonright_{X}, \tilde{u}=u_{1}+u_{2} \in$ $J_{s+\frac{1}{2}} L_{\text {loc }}^{2}(\tilde{X}, H)$. Since $H$ varies freely in $\mathscr{R}$ we conclude the proof by observing that the singular support condition required by Definition 3.6 is easily furnished by Lemmas 8.1 and 8.2.

The proof of the main theorem on the conormal regularity for semi-linear diffractive mixed problems is now an equally easy consequence of Theorems $3.8,7.23$ and 8.3. We recall that $F$ is a $C^{\infty}$ characteristic surface satisfying $F \cap \bar{X}_{-} \cap \partial X=\emptyset$ and the pseudo-conormal space $J_{k} L_{\text {loc }}^{2}(X)$ is given by Definition 3.6.
Theorem 8.4. - Let $u \in L_{\mathrm{loc}}^{\infty}(X)$ be the solution of the semi-linear mixed problem:

$$
P u=f(x, u) \quad \text { in } \quad X, \quad u \upharpoonright_{\partial X}=0, \quad u \upharpoonright_{X_{-}}=u_{0}
$$

where $f \in C^{\infty}(\mathbb{C})$ and $u_{0} \in I_{k} L_{\text {loc }}^{2}\left(X_{-}, F\right)$. Then $u \in J_{k} L_{\text {loc }}^{2}(X)$.
Proof. - To apply the standard procedure based on the algebra property (Theorem 3.8 ) and the propagation property (Theorem 8.3) (see [32, 34] and references given there) we only need to check that

$$
P w=0, w \upharpoonright_{\partial X}=0, w \upharpoonright_{X_{-}}=u_{0} \Longrightarrow w \in J_{k} L_{\mathrm{loc}}^{2}(X)
$$

That however is easy now as Theorem 7.23 shows that $w \in J_{k} L_{\text {loc }}^{2}(\tilde{X}, H) \upharpoonright_{X}$, for any $H \in \mathscr{R}$, and the singular support statement follows from propagation of singularities for the diffractive mixed problem (cf. [35] and [14], Sect. 24.4).

The results presented in chapter 1 are easy consequences of Theorem 8.4 and the definitions.

## A. GLANCING HYPERSURFACES AND $b$-GEOMETRY

The purpose of this appendix is to present some refinements of the equivalence of glancing hypersurfaces (see chapter 2 for definitions and [24, 35] for a detailed discussion) and of the construction of solutions to the diffractive transport equations.

We will use the notation similar to that in chapter $4: \mathbb{R}_{+}^{n+1}=\{(x, y): x>$ $\left.0, y \in \mathbb{R}^{n}\right\}$ and denote the coordinates in $T^{*} \mathbb{R}^{n+1} \backslash 0$ and ${ }^{b} T^{*} \overline{\mathbb{R}}_{+}^{n+1} \backslash 0$ by $(x, y ; \xi, \eta)$ and $(x, y ; \lambda, \eta)$ respectively, so that $\jmath: T^{*} \overline{\mathbb{R}}_{+}^{n+1} \backslash 0 \rightarrow{ }^{b} T^{*} \overline{\mathbb{R}}_{+}^{n+1} \backslash 0$ takes the form $\jmath(x, y ; \xi, \eta)=(x, y ; x \xi, \eta)$.

Our starting point is the following theorem from [35]:
Proposition 1.1. - If $P$ and $Q$ in $T^{*} \overline{\mathbb{R}}_{+}^{n+1} \backslash 0$ are given by

$$
Q=\{x=0\} \quad P=\{p=0\}, p=\xi^{2}+2 a(x, y, \eta) \xi+b(x, y, \eta), a \in S_{\mathrm{hg}}^{1}, b \in S_{\mathrm{hg}}^{2}
$$

and are glancing at $m_{0}=(0 ;(0,1,0, \cdots, 0))$, then there exists a conic neighbourhood of $\jmath\left(m_{0}\right), \Gamma \subset{ }^{b} T^{*} \overline{\mathbb{R}}_{+}^{n+1} \backslash 0$ and ab-canonical transformation

$$
{ }^{b} \chi: \Gamma \longrightarrow{ }^{b} T^{*} \overline{\mathbb{R}}_{+}^{n+1} \backslash 0, \quad{ }^{b} \chi\left(\jmath\left(m_{0}\right)\right)=\jmath\left(m_{0}\right)
$$

such that

$$
{ }^{b} \chi\left(\left\{\lambda^{2}+\epsilon x^{3} \eta_{1}^{2}-x^{2} \eta_{1} \eta_{n}=0\right\} \cap \Gamma\right) \subset \overline{\jmath(P)}, \quad \epsilon=\operatorname{sgn}\left(\partial_{x} b(0,0)\right)
$$

We remark that ${ }^{b} \chi(\{x=0\} \cap \Gamma) \subset\{x=0\}$ is immediately satisfied and that a comparison with the general discussion in chapter 2 shows that $\epsilon=-1$ and $\epsilon=1$ in the diffractive and gliding cases respectively. It is also important to remember that ${ }^{b} \chi$ is essentially obtained by extending an appropriately chosen equivalence of glancing hypersurfaces (or rather its restriction to the boundary, $\chi_{\partial}$ ):

$$
\chi: \Gamma_{0} \longrightarrow T^{*} \mathbb{R}^{n+1} \backslash 0, \Gamma_{0} \subset T^{*} \mathbb{R}^{n+1} \backslash 0, \quad \jmath(\chi(m))={ }^{b} \chi(\jmath(m))
$$

For $\chi$ we immediately have $\chi^{*} p=a(x, y ; \xi, \eta)\left(\xi^{2}+\epsilon x \eta_{1}^{2}-\eta_{1} \eta_{n}\right)$, where $a$ is homogeneous of degree 0 and non-zero in $\Gamma_{0}$. The corresponding statement for ${ }^{b} \chi$ does not however follow immediately from Proposition 1.1. Nevertheless we have

Proposition 1.2. $-\quad$ Let $p \in S_{\mathrm{hg}}^{2,2}\left({ }^{b} T^{*} \overline{\mathbb{R}}_{+}^{n+1} \backslash 0\right)$ satisfy

$$
\left.x^{2} p\right|_{\lambda^{2}-x^{3} \eta_{1}^{2}-x^{2} \eta_{n} \eta_{1}=0}, \quad d \jmath^{*} p\left(m_{0}\right) \neq 0 .
$$

Then, for some conic neighbourhood of $m_{0}=(0,(0,1,0, \cdots, 0)), \Gamma_{1} \subset T^{*} \overline{\mathbb{R}}_{+}^{n+1} \backslash 0$, there exists a b-canonical transformation

$$
{ }^{b} \chi_{1}: \Gamma_{1} \longrightarrow{ }^{b} T^{*} \overline{\mathbb{R}}_{+}^{n+1} \backslash 0
$$

such that

$$
\begin{equation*}
{ }^{b} \chi_{1}^{*} p=c\left(\xi^{2}-x \eta_{1}^{2}-\eta_{n} \eta_{1}\right), \quad c \in C^{\infty}\left({ }^{b} T^{*} \mathbb{R}_{+}^{n+1} \backslash 0\right), \quad c \neq 0 \text { in } \Gamma_{1} . \tag{1.1}
\end{equation*}
$$

Proof. - By the Malgrange preparation theorem (see [14], Sect. 7.5) we can write

$$
\begin{equation*}
x^{2} p=r_{0}(x, y, \eta)+r_{1}(x, y, \eta) \lambda+\left(\lambda^{2}-x^{3} \eta_{1}-x^{2} \eta_{1} \eta_{n}\right) s(x, y, \lambda, \eta), \tag{1.2}
\end{equation*}
$$

for $(x, y ; \lambda, \eta)$ in some conic neighbourhood of $\jmath\left(m_{0}\right)$. On the other hand, since $p \in S_{h g}^{2,2}$,

$$
x^{2} p=x^{2} s_{2}(x, y, \lambda, \eta)+\lambda x s_{1}(x, y, \lambda, \eta)+\lambda^{2} s_{0}(x, y, \lambda, \eta) .
$$

The differentiation of (1.2) in $x$ and putting $x=\lambda=0$ shows that $\partial_{x} r_{0}(0, y, \eta)=$ $r_{0}(0, y \eta)=0$. The comparison of the two expressions for $x^{2} p$ gives $r_{1}(0, y, \eta) \lambda+$ $\lambda^{2} s(0, y, \lambda, \eta)=\lambda^{2} s_{0}(0, y, \eta, \lambda)$ and consequently $r_{1}(0, y, \eta) \equiv 0$. Since $r_{0}+r_{1} \lambda$ vanishes identically when $\lambda^{2}=x^{2}\left(x \eta_{1}^{2}+\eta_{1} \eta_{n}\right)$, it also follows that $r_{0}$ and $r_{1}$ vanish identically when $x \eta_{1}+\eta_{n} \geq 0\left(\eta_{1}>0\right.$ near $\left.\jmath\left(m_{0}\right)\right)$. The assumption on $p$ implies that $s\left(m_{0}\right) \neq 0$ and thus, in view of the above discussion, we can assume that sufficiently near $m_{0}$

$$
p=\xi^{2}-x \eta_{1}^{2}-\eta_{n} \eta_{1}+\rho_{0}(x, y, \eta)+\xi \rho_{1}(x, y, \eta)
$$

where $\rho_{0}$ and $\rho_{1}$ vanish identically in $\eta_{n}+x \eta_{1} \geq 0$.
To eliminate the last two terms in $p$, we construct a $b$-canonical transformation by the homotopy method. Thus we define

$$
p_{s}=\xi^{2}-x \eta_{1}^{2}-\eta_{n} \eta_{1}+s\left(\rho_{0}(x, y, \eta)+\xi \rho_{1}(x, y, \eta)\right)
$$

and we want to find $a_{s}$ and $b_{s}$ homogeneous of degree 0 in $(\lambda, \eta)$ and such that

$$
\left.\frac{d}{d s}\left(\exp H_{b_{s}}\right)^{*}\left(a_{s} p\right)\right)=0
$$

which is the same as

$$
\left\{\frac{d b_{s}}{d s}, a_{s} p_{s}\right\}+\frac{d}{d s}\left(a_{s} p_{s}\right)=0
$$

By introducing $\alpha_{s}=\log a_{s}$ (the solution $a_{s}$ is required to be strictly positive) this reduces to

$$
\begin{equation*}
\left\{\frac{d b_{s}}{d s}, p_{s}\right\}+\frac{d p_{s}}{d s}=\left(-\left\{\frac{d b_{s}}{d s}, \alpha_{s}\right\}-\frac{d \alpha_{s}}{d s}\right) p_{s} \tag{1.3}
\end{equation*}
$$

We easily see that, with $\eta_{1}=1$ for simplicity,

$$
\begin{aligned}
\jmath_{*} H_{p_{s}}= & 2 \xi\left(\frac{\partial}{\partial x}+\xi \frac{\partial}{\partial \lambda}\right)+x \frac{\partial}{\partial \lambda}-2 x \frac{\partial}{\partial y_{1}} \\
& -s x \frac{\partial \rho}{\partial x} \frac{\partial}{\partial \lambda}+s \rho_{1}\left(\frac{\partial}{\partial x}+\xi \frac{\partial}{\partial \lambda}\right)+\widetilde{H}_{-\eta_{1} \eta_{n}+s\left(\rho_{0}+\xi \rho_{1}\right)}
\end{aligned}
$$

where $\widetilde{H}$ is the tangential Hamilton vector field. From this we compute for $q_{s} \in$ $C^{\infty}\left({ }^{b} T^{*} \overline{\mathbb{R}}_{+}^{n+1} \backslash 0\right)$

$$
\begin{aligned}
-x\left\{q_{s}, p_{s}\right\}= & \left(\lambda \frac{\partial}{\partial x}+2\left(\eta_{n}-s \rho_{0}(0, y, \eta)\right) x \frac{\partial}{\partial \lambda}+s \rho_{1}(0, y, \eta)\left(x \frac{\partial}{\partial x}-\lambda \frac{\partial}{\partial \lambda}\right)\right) q_{s} \\
& +p_{s} x \frac{\partial q_{s}}{\partial \lambda}+x \widetilde{H}_{-\eta_{1} \eta_{n}+s\left(\rho_{0}+\xi \rho_{1}\right)} q_{s}+W_{s} q_{s}
\end{aligned}
$$

where $W_{s}$ is a vectorfield with coefficients vanishing to second order at $x=\lambda=0$. With $q_{s}=d b_{s} / d s$ in mind (see (1.3) we want to solve

$$
\begin{equation*}
-x\left\{q_{s}, p_{s}\right\}=x \frac{d p_{s}}{d s}+x \mu_{s} p_{s} \tag{1.4}
\end{equation*}
$$

for some smooth $\mu_{s}=\mu_{s}(x, y, \lambda, \eta)$ so that $q_{s}$ and $\mu_{s}$ vanish in $x \eta_{1}+\eta_{n} \geq 0$.
To solve (1.4) we first construct the Taylor series of $q_{s}$ and $\mu_{s}$ at $x=\lambda=0$ and for that we need the following

Lemma 1.3. - Let $H_{l}$ be the space of real homogeneous polynomials of degree $l$ in $x, \lambda$ and let $A(t)$ be defined as

$$
A(t)=\lambda \frac{\partial}{\partial x}+t x \frac{\partial}{\partial \lambda}+f(t)\left(\lambda \frac{\partial}{\partial \lambda}-x \frac{\partial}{\partial x}\right)
$$

with $f \in C^{\infty}(\mathbb{R} ; \mathbb{R}), f(t)=\mathcal{O}\left(t^{N}\right)$, for all $N$ as $t \rightarrow 0$.
Then for $t \neq 0$ and $k \in \mathbb{N}_{0}$
i) The linear transformation $\left.A(t)\right|_{H_{2 k+1}}$ is invertible, $\left\|\left(\left.A(t)\right|_{H_{2 k+1}}\right)^{-1}\right\|=\mathcal{O}\left(|t|^{-1}\right)$.
ii) Every $u \in H_{2 k}$ can be written as

$$
u=A(t) v(t)+c(t)\left(\lambda^{2}-t x^{2}-2 f(t) x \lambda\right)^{k}, c(t) \in \mathbb{R}
$$

$\left\|(d / d t)^{m} v\right\|=\mathfrak{O}\left(|t|^{-1-m}\right)\|u\|, \quad(d / d t)^{m} c=\mathfrak{O}\left(|t|^{-1-m}\right)\|u\|$, for any fixed norms $\|\bullet\|=\|\bullet\|_{l}$ on $H_{l}$.

Proof. - To keep notation simple we will only consider the case $t<0$. It is convenient to introduce the following change of variables

$$
F(t)\binom{x}{y}=\binom{\sqrt{|t|} x}{y}=\binom{x_{1}}{y_{1}}
$$

so that

$$
\begin{equation*}
\left(F(t)^{-1}\right)^{*} A(t) F(t)^{*}=\sqrt{|t|}\left[\lambda_{1} \frac{\partial}{\partial x_{1}}-x_{1} \frac{\partial}{\partial \lambda_{1}}+\frac{f(t)}{\sqrt{|t|}}\left(\lambda_{1} \frac{\partial}{\partial \lambda_{1}}-x_{1} \frac{\partial}{\partial x_{1}}\right)\right] \tag{1.5}
\end{equation*}
$$

Let $\langle\bullet, \bullet\rangle_{l}$ denote the inner product on $H_{l}$ obtained by restricting to $x^{2}+\lambda^{2}=1$ and taking the $L^{2}$ inner product. We introduce a $t$-dependent inner product on $H_{l}$ :

$$
\left\langle F(t)^{*} v, F(t)^{*} w\right\rangle_{t, l}=\langle v, w\rangle_{l} .
$$

It depends smoothly on $t$ for $t<0$ and satisfies

$$
\langle v, w\rangle_{l} / C \leq\langle v, w\rangle_{t, l} \leq C|t|^{-1}\langle v, w\rangle_{l}, \quad\left(\frac{d}{d t}\right)^{m}\langle v, v\rangle_{t, l}=\mathcal{O}\left(|t|^{-1-m}\right)\langle v, v\rangle_{l} .
$$

Thus, we only need to consider $A_{1}(t)$ given by (1.5). If $l=2 k+1,\left.A_{1}(t)\right|_{H_{l}}$ is invertible with the norm bounded by $|t|^{-1 / 2}(1+o(1))$. Consequently, $\left.A(t)\right|_{H_{l}}$ is invertible with the norm bounded by $|t|^{-1}(1+o(1))$. When $l=2 k,\left.A_{1}(t)\right|_{H_{l}}$ has a one dimensional kernel spanned by $\left(\lambda_{1}+x_{1}^{2}-2 f(t) / \sqrt{|t|} \lambda_{1} x_{1}\right)^{k}$. Thus we can take the inverse of $A_{1}(t)$ restricted to the orthogonal complement of the kernel with respect to $\langle\bullet, \bullet\rangle_{l}$. Translating back to $A(t)$ that gives ii) with $c=\left\langle u,\left(\lambda^{2}-t x^{2}-2 f(t) x \lambda\right)^{k}\right\rangle_{l, t}$.

Since $s d p_{s} / d s=x \rho_{0}+\lambda \rho_{1}$ we can solve (1.4) in Taylor series:

$$
\begin{align*}
& q_{s}(x, \lambda ; x, \eta) \sim \sum_{l=1}^{\infty} q_{s}^{(l)}(x, \lambda ; y, \eta), \quad q_{s}^{(l)}(\bullet ; y, \eta) \in H_{l}, q_{s}^{(l)} \equiv 0 \text { in } \eta_{n} \geq 0  \tag{1.6}\\
& \mu_{s}(x, \lambda ; x, \eta) \sim \sum_{l=1}^{\infty} \mu_{s}^{(l)}(x, \lambda ; y, \eta), \quad \mu_{s}^{(l)}(\bullet ; y, \eta) \in H_{l}, \mu_{s}^{(l)} \equiv 0 \text { in } \eta_{n} \geq 0 \tag{1.7}
\end{align*}
$$

In fact, we apply Lemma 1.3 with $t=\eta_{n}-s \rho_{0}(0, y, \eta)$ and $f(t)=s \rho_{1}(0, y, \eta) / 2$ (treating $\eta_{n}$ as the only variable and the remaining ones as parameters), set $q_{s}^{(0)}=0$ and

$$
q_{s}^{(1)}(x, \lambda ; y, \eta)=\left(2 A(t) \mid H_{1}\right)^{-1}\left(x \rho_{0}(0,0 ; y, \eta)+\lambda \rho_{1}(0,0 ; y, \eta)\right), \quad \mu_{s}^{0}=\partial_{\lambda} q_{s}^{(1)}
$$

We then continue solving

$$
\begin{aligned}
2 A(t) q_{s}^{(2 k)}= & \gamma_{s}^{(2 k)}+c_{s}^{(2 k)}\left(\lambda^{2}-t x^{2}-2 f(t) x \lambda\right)^{k} \\
& \left\langle\gamma_{s}^{(2 k)},\left(\lambda^{2}-t x^{2}-2 f(t) x \lambda\right)^{k}\right\rangle_{t, 2 k}=0 \\
2 A(t) q_{s}^{(2 k+1)}= & \gamma_{s}^{(2 k+1)}, \quad c_{s}^{(2 k+1)}=0
\end{aligned}
$$

where $\gamma_{s}^{(l)}$ are obtained from $q_{s}^{(m)}$ with $m<l$. The corresponding $\mu_{s}^{(l)}$ are given by

$$
\begin{equation*}
\mu_{s}^{(l)}=c_{s}^{(l)} x\left(\lambda^{2}-t x^{2}-2 f(t) x \lambda\right)^{\frac{l-2}{2}}+\partial_{\lambda} q_{s}^{(l+1)} \tag{1.8}
\end{equation*}
$$

The crucial observations are that

$$
\lambda^{2}-t x^{2}-2 f(t) x \lambda=x^{2} p_{s}+\mathcal{O}\left(|x|^{3}+|\lambda|^{3}\right)
$$

and that the last term in the right hand side of (1.8) does not contribute to $\gamma_{s}^{(\bullet)}$. The norm estimates in Lemma 1.3 show that $q_{s}^{(l)}, \mu_{s}^{(l)}$ are $C^{\infty}$ in $(y, \eta)$ and the vanishing in $\eta_{n} \geq 0$ is immediate. The Borel lemma now gives $q_{s}, \mu_{s}, C^{\infty}$ in $(x, y ; \lambda, \eta)$ and $s$, vanishing identically in $\eta_{n}+\eta_{1} x \geq 0$ and such that (1.6) holds. Thus,

$$
-x\left\{q_{s}, p_{s}\right\}=x \frac{d p_{s}}{d s}+x p_{s} \mu_{s}+\nu_{s}
$$

where $\nu_{s}=\nu_{s}(x, y ; \lambda, \eta)$ vanishes in $\eta_{n}+\eta_{1} x$ and to infinite order at $x=\lambda=0$. But then $\nu_{s} /\left(x^{2} p_{s}\right) \in C^{\infty}\left({ }^{b} T^{*} \overline{\mathbb{R}}_{+}^{n+1} \backslash 0\right)$ and we have solved (1.4) with $\mu_{s}$ replaced by $\mu_{s}+x\left(\nu_{s} /\left(x^{2} p_{s}\right)\right)$.
Going back to (1.3) we now solve for $b_{s}$ and $\alpha_{s}$

$$
\begin{aligned}
& \frac{d b_{s}}{d s}=q_{s}, \quad b_{0}=0 \\
& \frac{d \alpha_{s}}{d s}-\left\{q_{s}, \alpha_{s}\right\}=-\mu_{s}, \quad \alpha_{0}=0
\end{aligned}
$$

Putting $a=\exp \left(\alpha_{1}\right)$ and ${ }^{b} \chi_{1}=\exp \left(H_{b_{1}}\right)$ we obtain the desired $b$-canonical transformation:

$$
{ }^{b} \chi_{1}^{*}(a p)=\xi^{2}-x \eta_{1}^{2}-\eta_{n} \eta_{1} .
$$

Remark 1.4. - Any b-canonical transformation ${ }^{b} \chi$ satisfies

$$
{ }^{b} \chi^{*} x=a x, \quad{ }^{b} \chi^{*} \lambda=b \lambda+c x, \quad a, b \neq 0
$$

and thus induces a canonical transformation on $T^{*} \mathbb{R}^{n} \backslash 0$ (or a conic subset of it):

$$
\chi_{\partial}(y ; \eta)=\left(y^{\prime}, \eta^{\prime}\right) \Longleftrightarrow{ }^{b} \chi(0, y ; 0, \eta)=\left(0, y^{\prime} ; 0, \eta^{\prime}\right)
$$

Since the construction in the proof of Proposition 1.2 gave $q_{s}(0, y ; 0, \eta) \equiv 0$ and consequently $b_{s}(0, y ; 0, \eta) \equiv 0,{ }^{b} \chi_{1}$ satisfying (1.1) can be chosen so that

$$
\left(\chi_{1}\right)_{\partial}=\left.I d\right|_{\Gamma_{0}}, \quad \Gamma_{0}=\left\{(y, \eta) \in T^{*} \mathbb{R}^{n} \backslash 0:(0, y ; 0, \eta) \in \Gamma_{1}\right\}
$$

This observation is very convenient in applications presented in chapter 7.
Combining the two propositions we obtain:

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Theorem 1.5. - The $b$-canonical transformation ${ }^{b} \chi$ in Proposition 1.1 can be chosen so that

$$
{ }^{b} \chi^{*} p=c\left(\xi^{2}-x \eta_{1}^{2}-\eta_{1} \eta_{n}\right)
$$

with $c \in S_{\mathrm{hg}}^{0}\left({ }^{b} T^{*} \overline{\mathbb{R}}_{+}^{n+1} \backslash 0\right), c \neq 0$ in $\Gamma$.
We now want to consider transport equations and again we start by recalling some facts already contained, in a slightly different form, in Sect. 4.4 of [35]:

Lemma 1.6. - If $\Gamma_{2} \subset \mathbb{R}_{x, y}^{n+1} \times \mathbb{R}_{\eta}^{n}$ is a conic neighbourhood of $m_{0}=(0 ;(1,0, \cdots, 0))$ (with respect to the $\mathbb{R}_{+}$action on the last $n$ variables) and $A, B \in S^{0}\left(\Gamma_{2}\right)$ then there exist a conic neighbourhood of $m_{0}, \Gamma_{3} \subset \Gamma_{2}$ and $g, h, e_{1}, e_{2} \in S^{0}\left(\Gamma_{2}\right)$ such that in $\Gamma_{3}$ :

$$
\left(\begin{array}{cc}
-\partial_{y_{n}}-\left(2+\eta_{1}^{-1} \eta_{n}\right) \partial_{y_{1}} & 1+2\left(x+\eta_{1}^{-1} \eta_{n}\right) \partial_{x} \\
2 \partial_{x} & -\partial_{y_{n}}-\left(2+\eta_{1}^{-1} \eta_{n}\right) \partial_{y_{1}}
\end{array}\right)\binom{g}{h}=\binom{A}{B}+\binom{e_{1}}{e_{2}}
$$

and the solution satisfies

$$
g\left(m_{0}\right)=1,\left.h\right|_{x=0} \equiv 0, e_{i} \equiv 0 \text { in } x \eta_{1}+\eta_{n} \geq 0, e_{i}=\mathfrak{O}\left(x^{N}\right), N \in \mathbb{N}, x \rightarrow 0
$$

This lemma allows to solve the transport equations when the operator is the model one: $p=\xi^{2}-x \eta_{1}^{2}-\eta_{1} \eta_{n}$. Here we are interested in the following slight refinement of the standard procedure:

Proposition 1.7. - If $p=\xi^{2}-x \eta_{1}^{2}-\eta_{1} \eta_{n}$ and $r \in S^{1,1}\left({ }^{b} T^{*} \overline{\mathbb{R}}_{+}^{n+1} \backslash 0\right)$ then there exist $a$ conic neighbourhood of $(0 ;(0,1,0, \cdots, 0)), \Gamma \subset{ }^{b} T^{*} \overline{\mathbb{R}}_{+}^{n+1} \backslash 0$ and $a \in S^{0}\left({ }^{b} T^{*} \overline{\mathbb{R}}_{+}^{n+1} \backslash 0\right)$, $b \in S^{-1}\left({ }^{b} T^{*} \overline{\mathbb{R}}_{+}^{n+1} \backslash 0\right)$ such that

$$
H_{p} a=r+p b \text { in } J^{-1}(\Gamma)
$$

where we identified $a, b, r$ with their pullbacks under $\jmath: T^{*} \overline{\mathbb{R}}_{+}^{n+1} \backslash 0 \rightarrow{ }^{b} T^{*} \overline{\mathbb{R}}_{+}^{n+1} \backslash 0$.
Proof. - Arguing as in the beginning of the proof of Proposition 1.2 (with the argument applied to $x r)$, we see that for $(x, y, x \xi, \eta) \in \Gamma$,

$$
r=r_{1}(x, y, \eta)+\xi r_{0}(x, y, \eta)+p b^{\sharp}(x, y ; x \xi, \eta),
$$

where $r_{i} \in S^{i}\left(\Gamma_{2}\right)$ with $\Gamma_{2} \subset \mathbb{R}^{n+1} \times \mathbb{R}^{n}$, a conic neighbourhood of $(0 ;(0,1,0, \cdots, 0)$ and $b^{\sharp} \in S^{-1}(\Gamma)$.
Let us now apply Lemma 1.6 with $A=\eta_{1}^{-1} r_{1}$ and $B=r_{0}$ to obtain $g$ and $h$. A simple computation shows that

$$
H_{p}\left(g+\xi \eta_{1}^{-1} h\right)=r_{1}+\xi r_{0}+2 p \eta_{1}^{-1} \partial_{x} h+e_{1}+\xi \eta_{1}^{-1} e_{2}
$$

Since $\left.h\right|_{x=0}=0$ we can take $a=g+\lambda h / x$, while $b=b^{\sharp}+2 \partial_{x} h+x^{2}\left[\left(e_{1}+\right.\right.$ $\left.\lambda \eta_{1}^{-1} e_{2} / x\right) /\left(x^{2} p\right)$ ], where the last term is smooth as $e_{i}$ vanish identically in $\eta_{1} x+\eta_{n} \geq$ 0 and to infinite order at $x=0$.

## B. $b$-SOBOLEV SPACES

In this appendix we recall the definition of the $b$-Sobolev spaces on a manifold with corners and outline the proof of some of their more useful properties. We start with the case of a manifold with boundary, for a discussion of the characteristic operators used, see [25, 28] and also [14], Sect. 18.3.

Let $M$ be a manifold with boundary and let $\nu$ be a measure on $M$. Then we define

$$
\begin{equation*}
H_{(s), \nu}^{b}(M)=\left\{u ; A u \in L_{\nu}^{2}(M) \text { for all } A \in \Psi_{b}^{s}(M)\right\} \tag{2.1}
\end{equation*}
$$

When the measure $\nu$ is not specified and $M$ is locally described by a subset of $[0, \infty)_{r} \times \mathbb{R}_{y}^{n-1}$, we take $\nu=d y d r / r$ and denote the corresponding $L_{\nu}^{2}$ by $L_{b}^{2}(M)$. A logarithmic change of variables $t=\log r$ induces an isomorphism

$$
\begin{equation*}
H_{s}^{b}\left([0, \infty) \times \mathbb{R}^{n-1}\right) \longleftrightarrow H_{(s)}\left(\mathbb{R}^{n}\right) \tag{2.2}
\end{equation*}
$$

and thus we easily conclude that $H_{(s), \nu}^{b}(M)$ is an interpolation space for any smooth measure $\nu$.


Figure 2.1. An example of the resolution $M_{1} \xrightarrow{\gamma} M$
More generally if $M$ is a manifold with corners, $H_{(s), \nu}^{b}(M)$ can be defined by (2.1) where $\Psi_{b}^{s}(M)$ is the space of totally characteristic operators on a manifold with corners [31] (see also [28]) or, more directly, by (2.2) and a multiple logarithmic

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change of variables - see [29]. As for manifolds with boundaries (see [14], Appendix B) more options for defining Sobolev spaces are available. To describe them let us assume that $M$ has a codimension two corner and near it, is described by

$$
Z_{2}=[0, \infty) \times[0, \infty) \times \mathbb{R}^{n-2}, \quad \text { with coordinates }(r, x, y) \in Z_{2}
$$

For that manifold we define

$$
\begin{aligned}
& \bar{H}_{(s)}^{b}\left(Z_{2}\right)=\left\{u ; \text { there exists } \tilde{u} \in H_{(s)}^{b}\left([0, \infty) \times \mathbb{R}^{n-1}\right) \text { such that } \tilde{u} \upharpoonright_{x>0}=u\right\} \\
& H_{(s, m)}^{b}\left(Z_{2}\right)=\left\{u: U \in H_{(s, m)}\left(\mathbb{R} \times[0, \infty) \times \mathbb{R}^{n-2}\right) \text { where } U(t, x, y)=u(\log r, x, y)\right\}
\end{aligned}
$$

where $H_{(s, m)}$ is defined in Appendix B of [14]. We easily have the following adaptation of Theorem B.2.9 to this setting

Proposition 2.1. - If $P \in \operatorname{Diff}_{b}^{p}\left([0, \infty) \times \mathbb{R}^{n-1}\right)$ and $\{x=0\} \subset Z_{2} \times \mathbb{R}^{n-2}$ is noncharacteristic for $P$, i.e. ${ }^{b} \sigma_{p}(P) \upharpoonright^{b_{N^{*}}\{x=0\}}$ does not vanish, then

$$
u \in H_{\left(s_{1}, m_{1}\right)}^{b}\left(Z_{2}\right) \text { and } P u \in H_{\left(s_{2}-p, m_{2}\right)}^{b}\left(Z_{2}\right) \Longrightarrow u \in H_{(s, m)}^{b}\left(Z_{2}\right)
$$

for $s+m \leq s_{i}+m_{i}, s \leq s_{2}$.
If $M$ is a manifold with corners and

$$
\begin{align*}
\mathscr{R}=R_{0} \sqcup & \cdots \sqcup R_{m}= \\
& \left\{\bigcup_{i} R_{i} \backslash\left(\bigcup_{i \neq j} R_{i} \cap R_{j}\right), R_{i} \cap R_{j} \backslash\left(\bigcup_{k} R_{i} \cap R_{j} \cap R_{k}\right), \cdots\right\} \tag{2.3}
\end{align*}
$$

is a variety of cleanly intersecting smooth submanifolds of $M, R_{0}=\partial M$, we define the conormal space $I_{k} L_{\nu}^{2}(M, \mathscr{R})$ by (1.1) using all vector fields tangent to $\mathcal{R}$ for ${ }^{\circ}$. By succesively blowing up all intersections (see [28]) and then the submanifolds $R_{i}, i>0$, we obtain the resolution $M_{1}$ :

$$
M_{1} \xrightarrow{\gamma} M,
$$

where $M_{1}$ is a manifold with corners - see Fig. 2.1 for an example. We then clearly see that

$$
I_{k} L_{\nu}^{2}(M, \mathscr{R}) \stackrel{\gamma_{*}}{\longleftrightarrow} H_{(k), \nu_{\gamma}}^{b}\left(M_{1}\right), \quad \gamma_{*} \nu_{\gamma}=\nu,
$$

and thus we obtain
Proposition 2.2. - If $M$ is a manifold with corners and $\mathcal{R}$ is given by (2.3) then $I_{k} L_{\nu}^{2}(M, \mathscr{R})$ are interpolation spaces in $k$.

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[^1]:    ${ }^{1}$ See chapter 2 ; we will refer to this artificial surface only once, in the last part of the proof of Theorem 7.6.

