

ON A B.V.P. IN A QUADRANT WITH  
NON-LOCAL BOUNDARY CONDITIONS

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**Abstract:** We consider a boundary value problem (b.v.p.) in a quadrant with absorbing boundary value conditions introduced by Engquist and Majda [3]. We prove the uniqueness of the solution in the Schwartz space of distributions.

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### 1. Introduction

The aim of this paper is to prove that the non-local boundary value conditions, introduced by Engquist and Majda [3], [4] lead to the well-posed boundary value problem for the Helmholtz equation in a quadrant.

These conditions appeared as some pseudo-differential operators that turn the boundary of the half-plane into transparent one for some traveling waves. After the Fourier transform with respect to  $t$ , this operator has the form

$$Bu(x, y) = \left( \frac{d}{dx} - i\sqrt{\omega^2 + \frac{d^2}{dy^2}} \right) u(x, y), \quad (x, y) \in \mathbb{R}^2, \quad \omega > 0, \quad (1)$$

where for  $v \in \mathcal{S}(\mathbb{R})$  (the Schwartz space of the smooth rapidly decreasing functions)

$$\sqrt{\omega^2 + \frac{d^2}{dy^2}} v(y) := \frac{1}{2\pi} \int_{\mathbb{R}} e^{-iy\eta} \sqrt{\omega^2 - \eta^2} \hat{v}(\eta) d\eta, \quad (2)$$

with

$$\hat{v}(\eta) = \int_{\mathbb{R}} e^{i\eta y} v(y) dy. \quad (3)$$

Let  $\mathcal{S}(\mathbb{R}^n)$  and  $\mathcal{S}'(\mathbb{R}^n)$  be the Schwartz spaces of the the smooth rapidly decreasing functions and the tempered distributions respectively. Then  $B$  is well-defined operator from  $\mathcal{S}(\mathbb{R}^2)$  to  $\mathcal{S}'(\mathbb{R}^2)$ . By continuity it admits the extension to the operator from  $\mathcal{S}'(\mathbb{R}^2)$  to  $\mathcal{S}'(\mathbb{R}^2)$ . This extension has the following remarkable property: it annihilates the travelling to the left plane waves on the line  $x = 0$ . These waves have the form  $u = e^{i\omega t} e^{-i\omega \vec{n} \cdot (x, y)}$ , with  $\vec{n} = (\cos \phi, \sin \phi)$ ,  $\phi \in (\pi/2, 3\pi/2)$ .

The approximations of this operator by the differential linear operators are used as approximate absorbing boundary conditions in numerical computations with artificial boundaries, introduced to limit a computational domain, (see [5] for an overview). These approximations may lead the boundary value problem in a half-plane to either well-posed or ill-posed problem according to the criterion introduced in [3] and [10].

There are several papers where the well-posed problems are studied for these boundary conditions, for example [10] for the half-plane, [1] for a non-stationary problem in a quadrant. We prove the uniqueness of the stationary problem for the *exact* non-local boundary value conditions in a quadrant.

Namely, we prove that in a quadrant the boundary value problem for the Helmholtz equation is well-posed in the following sense: there exist no non-trivial solutions of the Helmholtz equation for  $x \in K := \{(x_1, x_2) \in \mathbb{R}^2 \mid x_1 > 0, x_2 > 0\}$

$$(\Delta + \omega^2)u(x) = 0, \quad \omega > 0, \quad (4)$$

with the boundary conditions of the type (1) in the Schwartz space of *distributions*. In this case we need to give an appropriate definition of the operator (1) because this operator is defined *a priori* on the functions given on  $\mathbb{R}^2$ . We do it by means of the Paley-Wiener theory using [9].

Note that the function

$$u(x_1, x_2) := \sin\left(\frac{\omega}{\sqrt{2}} x_1\right) \sin\left(\frac{\omega}{\sqrt{2}} x_2\right),$$

is a solution to the boundary value problem for the Helmholtz equation (4) with the Dirichlet boundary condition in the quadrant  $K$ . This function belongs to the space  $\mathbb{C}^\infty(\overline{K})$  of the smooth functions in  $\overline{K}$  that are bounded with all their derivatives. Moreover,  $u(x)$  is equal to 0 on the lattice  $L := \left\{ (x_1, x_2) \in \mathbb{R}^2 \mid x_1 = \frac{k\pi\sqrt{2}}{\omega}, k \in \mathbb{Z}, x_2 \in \mathbb{R} \right\} \cup \left\{ (x_1, x_2) \in \mathbb{R}^2 \mid x_2 = \frac{k\pi\sqrt{2}}{\omega}, k \in \mathbb{Z}, x_1 \in \mathbb{R} \right\}$ . In particular, this is the reason why the Dirichlet boundary conditions are not suitable for computing simulations of the wave propagation problems (see [6], [2]). Here we prove that the same is not possible in a quadrant if the boundary operators are the modified operator of the type (1).

## 2. Fourier-Laplace Transform

Let us introduce some notations. Let  $\mathbb{R}^+ := \{x \in \mathbb{R} \mid x > 0\}$ ;  $K_+^n := \{x \in \mathbb{R}^n \mid x_i > 0, i = 1, \dots, n\}$ ,  $n = 1, 2$ ;  $K_+^1 := \mathbb{R}^+$ ,  $K_+^2 := K$ . We use the Fourier transform  $\mathbb{F} : \mathbb{S}(\mathbb{R}^n) \rightarrow \mathbb{S}(\mathbb{R}^n)$  in the form

$$\hat{f}(\xi) \equiv \mathbb{F}_{x \rightarrow \xi}[f](\xi) := \int_{\mathbb{R}^n} e^{i\langle \xi, x \rangle} f(x) dx, \quad \xi \in \mathbb{R}^n. \quad (5)$$

This transform is extended by the continuity to the Fourier transform of the space  $\mathbb{S}'(\mathbb{R}^n)$ . By  $\mathbb{S}'(\overline{K_+^n})$  we denote the space of the tempered distributions with the support in  $\overline{K_+^n}$ , and by  $\mathbb{S}'(K_+^n)$  we denote the space of its restrictions to the open region  $K_+^n$ :  $\mathbb{S}'(K_+^n) := \mathbb{S}'(\mathbb{R}^2)|_{K_+^n}$ . The inverse transform reads as

$$\mathbb{F}_{\xi \rightarrow x}^{-1}[\hat{f}](x) := \frac{1}{(2\pi)^n} \int_{\mathbb{R}^n} e^{-i\langle \xi, x \rangle} \hat{f}(\xi) d\xi, \quad x \in \mathbb{R}^n. \quad (6)$$

For an open region  $A$  in  $\mathbb{C}^n$  the symbol  $\mathcal{H}(A)$  denotes the set of all analytic functions in  $A$ , and  $\mathbb{C}^+ := \{z \in \mathbb{C} \mid \text{Im } z > 0\}$ . We will use the complex Fourier-Laplace transform and a version of the Paley-Wiener Theorem.

**Theorem 2.1.** (Paley-Wiener, see [8]) *i) Let  $f \in \mathbb{S}'(\overline{K_+^n})$ . Then  $\hat{f}(\xi)$  has an unique analytic extension,  $\tilde{f}(z)$ , to the domain  $\mathbb{C}K_+^n := \mathbb{R}^n \oplus iK_+^n$ ,  $z := (z_1, \dots, z_n) = (\xi_1 + i\tau_1, \dots, \xi_n + i\tau_n)$ ,  $\xi_i, \tau_i \in \mathbb{R}^+$ ,  $i = 1, \dots, n$ , such that*

$$\left| \tilde{f}(z) \right| \leq C(1 + |z|)^\mu \text{dist } \{z, \mathbb{R}^n\}, \quad z \in \mathbb{C}K_+^n, \quad (a)$$

for some constants  $\mu \in \mathbb{R}$ ,  $C > 0$  with  $\text{dist } \{z, \mathbb{R}^n\} := \min_{i=1, \dots, n} \{\text{Im } z_i\}$ , for  $z \in \mathbb{C}K_+^n$  and

$$\tilde{f}(\xi_1 + i\tau_1, \dots, \xi_n + i\tau_n) \xrightarrow{\tau_{1, \dots, n} \rightarrow 0^+} \hat{f}(\xi_1, \dots, \xi_n) \quad (b)$$

in  $\mathcal{S}'(\mathbb{R}^n)$ .

ii) Conversely, if  $\tilde{f}(z)$  is analytic function on  $\mathbb{C}K_+^n$  and (a), (b) hold, then there exists

$$\hat{f}(\xi) = \lim_{\text{Im } z \rightarrow 0^+} \tilde{f}(z),$$

in  $\mathcal{S}'(\mathbb{R}^n)$  and it is the Fourier transform of a distribution  $f(x)$  belonging to  $\mathcal{S}'(\overline{K_+^n})$ .

If  $\hat{f}(\xi) = \lim_{\text{Im } z \rightarrow 0^+} \tilde{f}(z)$ , then we write

$$\tilde{f}(z) := \mathbb{F}_{x \rightarrow z}[f(x)], \quad f(x) = \mathbb{F}_{z \rightarrow x}^{-1}[\tilde{f}(z)]$$

and  $\mathbb{F}$  is the complex Fourier-Laplace transform for  $f \in \mathcal{S}'(\overline{K_+^n})$ . In the next section we introduce the absorbing boundary value conditions.

### 3. Absorbing Boundary Value Conditions and Main Result

We modify the definition of the absorbing boundary value conditions proposed in [3]. For this we use the method [9] of the representation of boundary value conditions via Cauchy data.

We are going to define the absorbing boundary value conditions for functions defined on  $K$ , using Theorem 2.1. Namely, we define the absorbing operator as follows.

Let  $u(x) \in \mathcal{S}'(K)$  be a solution of Helmholtz equation (4). Then  $u \in \mathbb{C}^\infty(K)$  by ellipticity, and its traces on the rays  $x_i = \epsilon$ ,  $\epsilon > 0$ ,  $i = 1, 2$ , parallel to the side of the quadrant  $K$  are well defined. It is easy to demonstrate that these traces belong to  $\mathcal{S}'(\mathbb{R}^+)$ . The same is valid for all derivatives of  $u$  on these rays. We use the following theorem proved in [9].

**Theorem 3.1.** (see Proposition 2.3, Proposition 2.4 [9]) *i) There exist the Cauchy data of the solution  $u$ :*

$$\left\{ \begin{array}{l} u(x_1, 0+) := u(x_1, x_2)|_{x_2=0+} = \lim_{\epsilon \rightarrow 0^+} u(x_1, \epsilon), \quad x_1 > 0, \\ u(0+, x_2) := u(x_1, x_2)|_{x_1=0+} = \lim_{\epsilon \rightarrow 0^+} u(\epsilon, x_2), \quad x_2 > 0, \\ \partial_{x_2} u(x_1, 0+) := \partial_{x_2} u(x_1, x_2)|_{x_2=0+} = \lim_{\epsilon \rightarrow 0^+} (\partial_{x_2} u(x_1, x_2)|_{x_2=\epsilon}), \quad x_1 > 0, \\ \partial_{x_1} u(0+, x_2) := \partial_{x_1} u(x_1, x_2)|_{x_1=0+} = \lim_{\epsilon \rightarrow 0^+} (\partial_{x_1} u(x_1, x_2)|_{x_1=\epsilon}), \quad x_2 > 0, \end{array} \right. \quad (7)$$

where the limits are understood in the sense  $\mathcal{S}'(\mathbb{R}^+)$ .

ii) There exists the canonical extension,  $u_0(x) := u_0(x_1, x_2) \in \mathcal{S}'(\overline{K})$  such that

$$\gamma_0(x) := (\Delta + \omega^2)u_0(x), \quad x \in \mathbb{R}^2, \quad (8)$$

and has the form for  $x \in \mathbb{R}^2$

$$\gamma_0(x) = v_1^1(x_1)\delta(x_2) + v_2^1(x_2)\delta(x_1) + v_1^0(x_1)\delta'(x_2) + v_2^1(x_2)\delta'(x_2), \quad (9)$$

where  $v_l^\beta(x_l) \in \mathcal{S}'(\overline{\mathbb{R}^+})$ ,  $l = 1, 2$ ,  $\beta = 0, 1$  are the extensions by zero of  $u_l^{(\beta)}(x)|_{x_l}$ , i.e.

$$\begin{aligned} v_1^0(x_1) &= u(x_1, 0+), \quad v_1^1(x_1) = \partial_{x_2}u(x_1, 0+), \quad x_1 > 0, \\ v_2^0(x_2) &= u(0+, x_2), \quad v_2^1(x_2) = \partial_{x_1}u(0+, x_2), \quad x_2 > 0. \end{aligned} \quad (10)$$

Note that the Cauchy data (7) are defined only on  $\mathbb{R}^+$ , and its extensions by 0,  $v_l^\beta(x_l)$ ,  $l = 1, 2$ ,  $\beta = 0, 1$ , are defined on  $\mathbb{R}$ .

Let  $\omega > 0$ . We denote  $\Gamma_1 := \{(x_1, 0) \mid x_1 \geq 0\}$ ,  $\Gamma_2 := \{(0, x_2) \mid x_2 \geq 0\}$ . Let  $\mathbb{L} := (-\infty, -\omega] \cup [\omega, \infty)$ .

**Definition 3.2.** We define  $\sqrt{\omega^2 - z^2} : \mathbb{C} \setminus \mathbb{L} \rightarrow \mathbb{C}$  as the analytic branch of  $\sqrt{\cdot}$ , such that  $\operatorname{Re} \sqrt{\omega^2 - z^2} > 0$ , for all  $z \notin \mathbb{L}$ .

Now we are able to define  $Bu|_{\Gamma_l}$ ,  $l = 1, 2$ .

**Definition 3.3.** Let  $v(x) \in \mathcal{S}'(\overline{\mathbb{R}^+})$ . We define the pseudo-differential operator  $\sqrt{\omega^2 + \frac{\partial^2}{\partial x^2}} : \mathcal{S}'(\overline{\mathbb{R}^+}) \rightarrow \mathcal{S}'(\overline{\mathbb{R}^+})$  as

$$\sqrt{\omega^2 + \frac{\partial^2}{\partial x^2}} [v(x)] := \mathbb{F}_{z \rightarrow x}^{-1} [\sqrt{\omega^2 - z^2} \tilde{v}(z)], \quad \operatorname{Im} z > 0.$$

This definition is correct by Theorem 2.1. Now we define the non-local boundary value operators.

**Definition 3.4.** Let  $u(x_1, x_2)$  be a solution of the equation (4) and  $v_l^\beta$ , ( $l = 1, 2$ ,  $\beta = 1, 2$ ) be its extended Cauchy data. We define

$$B_l u(x_1, x_2)|_{\Gamma_l} := v_l^1(x_l) - i \sqrt{\omega^2 + \frac{\partial^2}{\partial x_l^2}} v_l^0(x_l), \quad l = 1, 2.$$

Note that for  $u \in \mathcal{C}^\infty(\overline{K})$  this definition coincides with  $Bu(x_1, x_2)|_{\Gamma_l}$ , where  $B$  is defined by (1).

In such a way we define the boundary operator  $(B_1, B_2)$  on the generalized solutions of the Helmholtz equation in quadrant, with the values in the space  $\mathcal{S}'(\overline{\mathbb{R}^+}) \times \mathcal{S}'(\overline{\mathbb{R}^+})$ .

Now let us consider the following homogeneous boundary value problem for the Helmholtz equation for  $u \in \mathbb{S}'(K)$ :

$$\begin{cases} (\Delta + \omega^2)u(x) = 0, & x \in K, \\ B_l u(x)|_{\Gamma_l} = 0, & l = 1, 2, \end{cases} \quad (11)$$

where  $B_l$ ,  $l = 1, 2$ , are defined by Definition 3.4. The goal is to prove the following main theorem.

**Theorem 3.5.** *Let  $\omega > 0$ . The problem (11) has only a trivial solution in  $\mathbb{S}'(K)$ .*

In the rest of the paper we prove this theorem.

#### 4. Reducing to an Algebraic Equation on the Riemann Surface

We are going to apply the complex Fourier transform  $\mathbb{F}_{x \rightarrow z}$  to (9). Let  $u \in \mathbb{S}'(K)$  be a solution to the problem (11). Then by Theorem 3.1 ii) there exists its extension  $u_0 \in \mathbb{S}'(\overline{K})$  such that (8) holds, moreover  $\text{supp } \gamma_0 \subset \partial K$ . Therefore, by Theorem 2.1,  $\tilde{u}_0(z)$  and  $\tilde{\gamma}_0(z)$  are holomorphic functions of two complex variables in the tube domain  $\mathbb{C}K$ . Hence, (8) implies

$$\begin{aligned} \mathbb{H}(z)u_0(z) &:= (-z_1^2 - z_2^2 + \omega^2)\tilde{u}_0(z_1, z_2) = \tilde{\gamma}_0(z) \\ &= \tilde{v}_1^1(z_1) + \tilde{v}_2^1(z_2) - iz_2\tilde{v}_1^0(z_1) - iz_1\tilde{v}_2^0(z_2), \\ &\quad (z_1, z_2) \in \mathbb{C}K. \end{aligned} \quad (12)$$

Here  $\tilde{v}_l^\beta(z_l) \in H(\mathbb{C}^+)$ ,  $l = 0, 1; \beta = 0, 1$  also by Theorem 2.1.

**Definition 4.1.** i)  $\mathcal{V}$  denotes the Riemann surface of the complex characteristics of the operator  $\Delta + \omega^2$

$$\mathcal{V} = \{z = (z_1, z_2) \in \mathbb{C}^2 \mid -z_1^2 - z_2^2 + \omega^2 = 0\}$$

ii)  $\mathcal{V}^* = \mathcal{V} \cap \mathbb{C}K$ .

**Lemma 4.1.** (Connection Equation) *Let  $u$  be a solution of (11). Then for the  $\gamma_0(x)$  from (9) the following identity holds*

$$\tilde{\gamma}_0(z) \equiv \tilde{v}_1^1(z_1) + \tilde{v}_2^1(z_2) - iz_2\tilde{v}_1^0(z_1) - iz_1\tilde{v}_2^0(z_2) = 0, \quad z \in \mathcal{V}^*. \quad (13)$$

*Proof.* The expression (12) holds for  $z \in \mathbb{C}K$ , in particular for  $z \in \mathcal{V}^*$ . But the left-hand side of (12) is zero for  $z \in \mathcal{V}^*$ , hence the right-hand side is also zero.  $\square$

We sum up our considerations in the following statement.

**Theorem 4.2.** *Let  $u(x)$  be a solution of (11). Then:*

i) *There exists  $u_0 \in S'(\overline{K})$  such that  $u_0$  satisfies (12).*

ii) *The function  $\tilde{\gamma}_0(z)$  fits (13). The solution  $u(x) = u_0|_K$  can be expressed in terms of  $\tilde{\gamma}_0(z)$  by the formula*

$$u_0(x) = \mathbb{F}_{z \rightarrow x}^{-1} \left[ \frac{\tilde{\gamma}_0(z)}{-z_1^2 - z_2^2 + \omega^2} \right], \quad x \in K, \quad (14)$$

where  $\mathbb{F}^{-1}$  is the inverse complex Fourier transformation.

Finally, let us apply the complex Fourier transform to the boundary conditions (11). By Definition 3.4 we obtain

$$\widetilde{B_l u_0(z)}|_{\Gamma_l} = \tilde{v}_l^1(z_l) - i\sqrt{\omega^2 - z_l^2} \tilde{v}_l^0(z_l) = 0, \quad z_l \in \mathbb{C}^+. \quad (15)$$

Now, eliminating  $\tilde{v}_1^1(z_1)$  and  $\tilde{v}_2^1(z_2)$  from (15) and substituting them to the connection equation (13), we get:

$$\left( \sqrt{\omega^2 - z_1^2} - z_2 \right) \tilde{v}_1^0(z_1) + \left( \sqrt{\omega^2 - z_2^2} - z_1 \right) \tilde{v}_2^0(z_2) = 0, \quad z \in \mathcal{V}^*. \quad (16)$$

We are going to prove that (16) has only trivial solution  $(\tilde{v}_1^0, \tilde{v}_2^0) \in \mathcal{H}(\mathbb{C}^+) \times \mathcal{H}(\mathbb{C}^+)$ . It means by (15) that  $\tilde{v}_l^1 \equiv 0$ ,  $l = 1, 2$ . Therefore  $\tilde{\gamma}_0 \equiv 0$  by (13) and  $u_0(x) \equiv 0$  by (14).

The idea is as follows. For  $z \in \mathcal{V}^*$   $z_1^2 + z_2^2 = \omega^2$ . It implies that  $\sqrt{\omega^2 - z_1^2} - z_2 = 0$  or  $\sqrt{\omega^2 - z_1^2} - z_2 = -2z_2$  depending on the sign of  $\sqrt{\cdot}$ . If  $\sqrt{\omega^2 - z_1^2} - z_2 = 0$  and  $\sqrt{\omega^2 - z_2^2} - z_1 \neq 0$  for  $(z_1, z_2)$  belonging to some open part of  $\mathcal{V}^*$ , then by analytic continuation  $\tilde{v}_2^0(z_2) \equiv 0$  in  $\mathbb{C}^+$ . Similarly  $\tilde{v}_1^0(z_1) \equiv 0$ ,  $\text{Im } z_1 > 0$ .

In the next sections we justify these arguments. Namely, we complement the equation (16) by two equations of automorphisms, which express that  $\tilde{v}_l^0(z_l)$  depends only on one variable. Then we determine the regions where the coefficients in (16) are zero. For this, it is convenient to introduce the uniformization parameter (see Section 6).

## 5. Automorphisms

Let

$$\mathcal{V}_l^+ := \left\{ z \in \mathcal{V} \mid z_l \in \mathbb{C}^+ \right\}, \quad l = 1, 2. \quad (17)$$

We define the natural lifting  $v_l^0(z)$  of  $\tilde{v}_l^0(z_l)$  on the Riemann surface  $\mathcal{V}$ :

$$v_l^0(z) := \tilde{v}_l^0(z_l), \quad z \in \mathcal{V}_l^+, \quad l = 1, 2. \quad (18)$$

Note that the functions  $v_l^0(z)$ ,  $l = 1, 2$ , are holomorphic in the region  $\mathcal{V}_l^+$ ,  $l = 1, 2$ , since the functions  $\tilde{v}_l^0(z_l)$ ,  $l = 1, 2$  are holomorphic in  $\mathbb{C}^+$ .

Equation (16) is an undetermined problem of a functional equation with two unknown functions. Nevertheless, it is a well-posed problem because of its structure: the function  $\tilde{v}_l^0$  depends only on  $z_l$ ,  $l = 1, 2$ . The latter can be expressed algebraically by the method [9]. Namely, the analytic functions  $v_l^0$  in  $\mathcal{V}_l^+$  are invariant with respect to the monodromy group of the covering  $p_l : \mathcal{V}_l^+ \rightarrow \mathbb{C}^+$ , where  $p_l z = z_l$ ,  $l = 1, 2$ . The monodromy group is the group of transformations  $\mathcal{V} \rightarrow \mathcal{V}$  that transpose the points  $z \in \mathcal{V}$  with the same coordinate  $z_l$ . These transformations transpose the “sheets of the covering”  $p_l$ , which are the roots of the characteristic equation  $H(z) = 0$ . Hence by Vieta theorem, we obtain the following formulas for the generators  $h_l : \mathcal{V} \rightarrow \mathcal{V}$  of the monodromy groups (automorphisms) of the covering  $p_l$

$$h_1(z_1, z_2) = (z_1, -z_2), \quad h_2(z_1, z_2) = (-z_1, z_2). \quad (19)$$

Therefore the invariance of  $v_l^0$  with respect to the  $h_l$  is written as

$$v_l^0(z) = v_l^0(h_l(z)), \quad z \in \mathcal{V}_l^+, \quad l = 1, 2. \quad (20)$$

We lift the equation (16) to  $\mathcal{V}^* = \mathcal{V}_1^+ \cap \mathcal{V}_2^+$ . Together with (20) we obtain the following system:

$$\begin{cases} \left( \sqrt{\omega^2 - z_1^2} - z_2 \right) v_1^0(z) + \left( \sqrt{\omega^2 - z_2^2} - z_1 \right) v_2^0(z) = 0, & z \in \mathcal{V}^*, \\ v_1^0(z) - v_1^0(h_1 z) = 0, & z \in \mathcal{V}_1^+, \\ v_2^0(z) - v_2^0(h_2 z) = 0, & z \in \mathcal{V}_2^+. \end{cases} \quad (21)$$

We have proved the following statement.

**Lemma 5.1.** *Let  $\tilde{v}_l^0(z_l) \in \mathcal{H}(\mathbb{C}^+)$ ,  $l = 1, 2$ , be a solution to (16). Then  $v_l^0(z)$ ,  $z \in \mathcal{V}_l^+$ ,  $l = 1, 2$ , is a solution to the system (21).*

We are going to prove that the system (21) admits only trivial solutions  $v_l^0 \in \mathcal{H}(\mathcal{V}_l^+)$ .



## 6. Universal Covering Surface

It is more convenient to use the uniformization parameter  $\mu$  for the Riemann surface  $\mathcal{V}$ . It is known that there exists a universal covering  $\check{\mathcal{V}}$  of the surface  $\mathcal{V}$  which is isomorphic to  $\mathbb{C}$  (see, e.g. [7], [11]). Let us denote  $\Pi : \check{\mathcal{V}} \rightarrow \mathcal{V}$  the natural projection. We introduce a parameter  $\mu \in \check{\mathcal{V}}$  such that

$$\Pi(\mu) = (z_1, z_2) = (-i\omega \sinh \mu, \omega \cosh \mu) \quad \mu \in \check{\mathcal{V}}. \quad (22)$$

Obviously,  $\Pi(\mu) \in \mathcal{V}$  for all  $\mu \in \check{\mathcal{V}}$ . Note that this transformation is periodic with the period  $2\pi i$ , hence for all  $z \in \mathcal{V}$ ,  $\Pi^{-1}(z) = \{\mu + 2i\pi k, k \in \mathbb{Z}\}$ .

**Definition 6.1.** Let

$$\check{\mathcal{V}}_l^+ := \left\{ \mu \in \check{\mathcal{V}} \mid \Pi(\mu) \in \mathcal{V}_l^+ \right\}, \quad l = 1, 2. \quad (23)$$

In the following lemma we give the explicit representation of  $\check{\mathcal{V}}_l^+$ .

**Lemma 6.1.** We have:

$$\check{\mathcal{V}}_1^+ := \check{\mathcal{V}}_{1R}^+ \cup \check{\mathcal{V}}_{1L}^+,$$

where

$$\begin{aligned} \check{\mathcal{V}}_{1R}^+ &= \left\{ \mu \in \check{\mathcal{V}} \mid \operatorname{Re} \mu > 0, \frac{\pi}{2} + 2k\pi < \operatorname{Im} \mu < \frac{3\pi}{2} + 2k\pi, k \in \mathbb{Z} \right\} \\ \check{\mathcal{V}}_{1L}^+ &= \left\{ \mu \in \check{\mathcal{V}} \mid \operatorname{Re} \mu < 0, -\frac{\pi}{2} + 2k\pi < \operatorname{Im} \mu < \frac{\pi}{2} + 2k\pi, k \in \mathbb{Z} \right\}. \end{aligned}$$

$$\check{\mathcal{V}}_2^+ = \check{\mathcal{V}}_{2R}^+ \cup \check{\mathcal{V}}_{2L}^+,$$

$$\check{\mathcal{V}}_{2R}^+ = \left\{ \mu \in \check{\mathcal{V}} \mid \operatorname{Re} \mu > 0, 2k\pi < \operatorname{Im} \mu < \pi + 2k\pi, k \in \mathbb{Z} \right\}$$

$$\check{\mathcal{V}}_{2L}^+ = \left\{ \mu \in \check{\mathcal{V}} \mid \operatorname{Re} \mu < 0, -\pi + 2k\pi < \operatorname{Im} \mu < 2k\pi, k \in \mathbb{Z} \right\}$$

(see Figure 1). and

$$\begin{aligned} \check{\mathcal{V}}^* &= \left\{ \mu \in \mathbb{C} \mid \operatorname{Re} \mu > 0, \frac{\pi}{2} + 2k\pi < \operatorname{Im} \mu < \pi + 2k\pi, k \in \mathbb{Z} \right\} \\ &\cup \left\{ \mu \in \mathbb{C} \mid \operatorname{Re} \mu < 0, -\frac{\pi}{2} + 2k\pi < \operatorname{Im} \mu < 2k\pi, k \in \mathbb{Z} \right\} \end{aligned}$$

(see Figure 2).

*Proof.* It follows from (22),  $\omega > 0$  and inequalities  $\operatorname{Im} z_l(\mu) > 0, l = 1, 2$ .  $\square$

Note that the regions  $\check{\mathcal{V}}_l^+, \check{\mathcal{V}}_{lR}^+, \check{\mathcal{V}}_{lL}^+, l = 1, 2$  are invariants with respect to the translation to  $2\pi i$ , because of the periodicity of the transformation  $\Pi(\mu)$ .

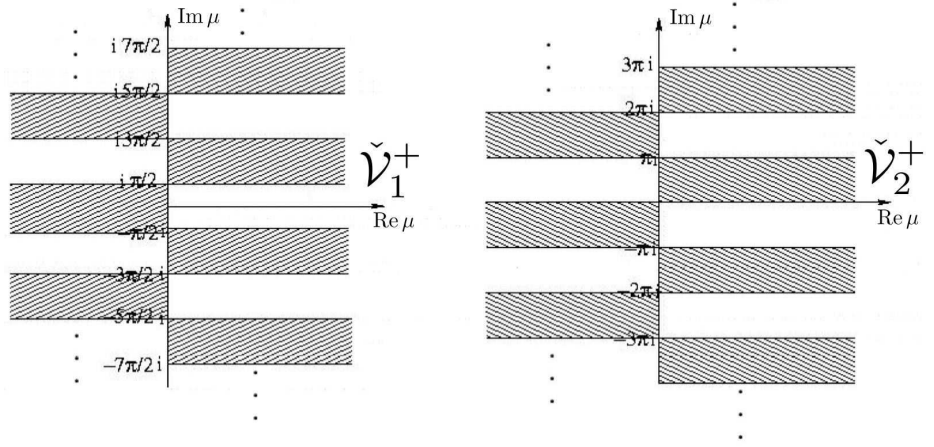


Figure 1: Regions  $\check{V}_1^+$  and  $\check{V}_2^+$

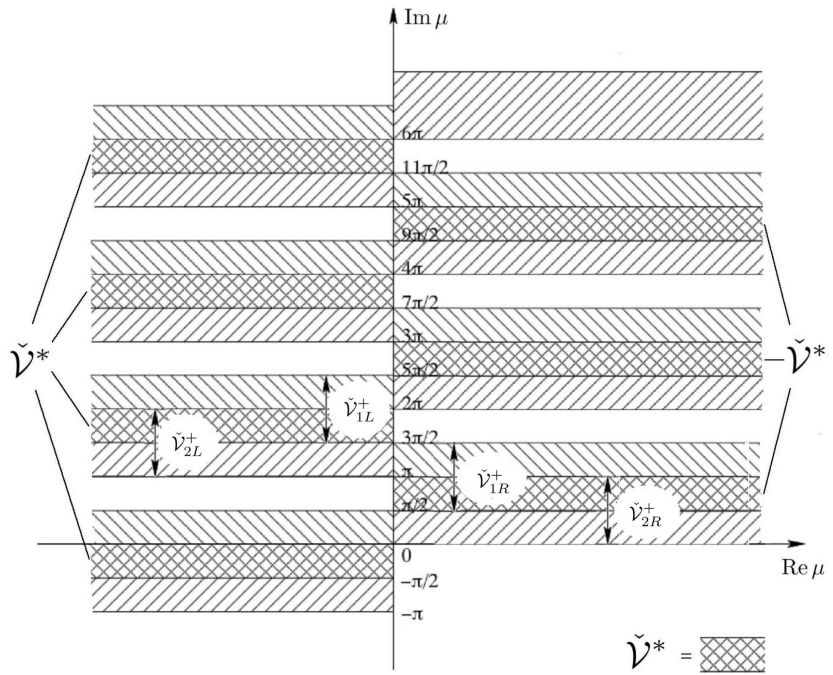


Figure 2: Description of the region  $\check{V}^*$

### 7. Lifting of the System (21) to the Universal Covering

**Lemma 7.1.** *There exist the automorphisms  $\check{h}_l : \check{V} \rightarrow \check{V}$ ,  $l = 1, 2$ , which are the lifting of the automorphisms  $h_l$ , i.e.*

$$\Pi \check{h}_l = h_l \Pi. \tag{24}$$

They act by formulas:

$$\check{h}_1(\mu) = -\mu + i\pi, \quad \check{h}_2(\mu) = -\mu, \quad (25)$$

i.e. the automorphism  $\check{h}_1$  is the symmetry of  $\check{\mathcal{V}}$  with respect to the point  $i\pi$ , and  $\check{h}_2$  is the symmetry with respect to the zero point.

*Proof.* It suffices to verify that the following diagram is commutative

$$\begin{array}{ccc} \check{\mathcal{V}} & \xrightarrow{\check{h}_l} & \check{\mathcal{V}} \\ \downarrow \Pi & & \downarrow \Pi \\ \mathcal{V} & \xrightarrow{h_l} & \mathcal{V} \end{array}$$

for  $h_l$  given by (25). This follows from (22) and (19). □

Further we lift the functions  $v_l^0(z)$  from (18) to  $\check{\mathcal{V}}$ . Let  $\Pi$  as before, then we define

$$\check{v}_l^0(\mu) = v_l^0(\Pi(\mu)), \quad \mu \in \check{\mathcal{V}}_l^+, \quad l = 1, 2. \quad (26)$$

Invariance (20), (24) and definition (26) imply that the functions  $\check{v}_l^0(\mu)$  are automorphic with respect to  $\check{h}_l$ , i.e.,

$$\check{v}_l^0(\mu) = \check{v}_l^0(\check{h}_l\mu), \quad \mu \in \check{\mathcal{V}}_l^+, \quad l = 1, 2, \quad (27)$$

where  $\check{h}_l$  are given by (25).

In the following we rewrite system (21) in terms of the variable  $\mu$ . Note that by the periodicity of all lifting it is sufficient to lift (21) to the principal region chosen arbitrarily. We denote subregion of  $\check{\mathcal{V}}$

$$\begin{aligned} \check{V}_{1R}^+ &:= \left\{ \operatorname{Re} \mu > 0, \operatorname{Im} \mu \in \left( \frac{\pi}{2}, \frac{3\pi}{2} \right) \right\}, \quad \check{V}_{2R}^+ := \left\{ \operatorname{Re} \mu > 0, \operatorname{Im} \mu \in (0, \pi) \right\}, \\ \check{V}_{1L}^+ &:= \left\{ \operatorname{Re} \mu < 0, \operatorname{Im} \mu \in \left( \frac{-\pi}{2}, \frac{\pi}{2} \right) \right\}, \quad \check{V}_{2L}^+ := \left\{ \operatorname{Re} \mu < 0, \operatorname{Im} \mu \in (-\pi, 0) \right\}. \end{aligned}$$

Also we denote

$$\begin{aligned} \check{V}_1^+ &:= \check{V}_{1R}^+ \cup \check{V}_{1L}^+, \quad \check{V}_2^+ := \check{V}_{2R}^+ \cup \check{V}_{2L}^+, \\ \check{V}_R^* &:= \check{V}_{1R}^+ \cap \check{V}_{2R}^+, \quad \check{V}_L^* := \check{V}_{1L}^+ \cap \check{V}_{2L}^+, \quad \check{V}^* = \check{V}_R^* \cup \check{V}_L^*, \end{aligned} \quad (28)$$

see Figure 3.

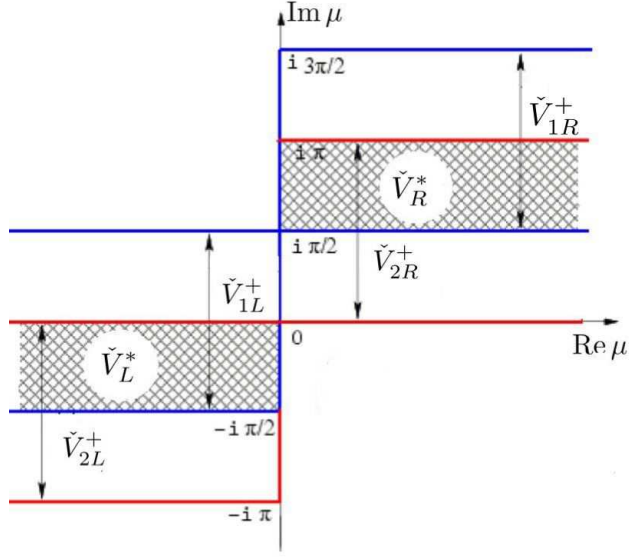


Figure 3: Principal region

It suffices to lift the system (21) to  $\check{V}_1^+ \cup \check{V}_2^+$ , where  $\check{V}_l^+$ ,  $l = 1, 2$  are defined by (28). Using (22), (26), (27) and (25) we obtain the equivalent system:

$$\left\{ \begin{array}{l} \left( \omega \sqrt{1 + \sinh^2 \mu} - \omega \cosh \mu \right) \check{v}_1^0(\mu) \\ \quad + \left( \omega \sqrt{1 - \cosh^2 \mu} + i \omega \sinh \mu \right) \check{v}_2^0(\mu) = 0, \quad \mu \in \check{V}^*, \\ \check{v}_1^0(-\mu + i\pi) = \check{v}_1^0(\mu), \quad \mu \in \check{V}_1^+, \\ \check{v}_2^0(-\mu) = \check{v}_2^0(\mu), \quad \mu \in \check{V}_2^+, \end{array} \right. \quad (29)$$

where  $\sqrt{\cdot}$  corresponds to the branch defined by Definition 3.2.

We prove that (29) has only trivial solution in the class  $\mathcal{H}(\check{V}_1^+) \times \mathcal{H}(\check{V}_2^+)$ . Dividing the first equation of (29) by  $\omega$  we obtain

$$\begin{aligned} \left( \sqrt{1 + \sinh^2 \mu} - \cosh \mu \right) \check{v}_1^0(\mu) \\ + \left( \sqrt{1 - \cosh^2 \mu} + i \sinh \mu \right) \check{v}_2^0(\mu) = 0, \quad \mu \in \check{V}^*. \end{aligned} \quad (30)$$

Now we analyze  $\sqrt{1 + \sinh^2 \mu}$ ,  $\mu \in \check{V}_1^+$  and  $\sqrt{1 - \cosh^2 \mu}$ ,  $\mu \in \check{V}_2^+$ .

**Lemma 7.2.** For  $\mu \in \check{V}^*$  we obtain

$$\sqrt{1 + \sinh^2 \mu} = \begin{cases} -\cosh \mu, & \mu \in \check{V}_R^*, \\ \cosh \mu, & \mu \in \check{V}_L^*. \end{cases} \quad (31)$$

$$\sqrt{1 - \cosh^2 \mu} = \begin{cases} -i \sinh \mu, & \mu \in \check{V}_R^*, \\ i \sinh \mu, & \mu \in \check{V}_L^*. \end{cases} \quad (32)$$

*Proof.* We consider the function

$$\sqrt{\omega^2 - z_1^2} : \mathbb{C}_{z_1}^+ \rightarrow \mathbb{C}.$$

This function is holomorphic in  $\mathbb{C}_{z_1}^+$ , hence its lifting by the formulas (22) is given by  $\sqrt{\omega^2 + \omega^2 \sinh^2 \mu}$  and is holomorphic in  $\check{V}_1^+$ . In general,

$$\sqrt{\omega^2 + \omega^2 \sinh^2 \mu} = \pm \omega \cosh \mu, \quad \mu \in \check{V}_1^+, \quad (33)$$

and we must choose the sign such that (33) corresponds to the branch fixed by Definition 3.2. It is easy to see that for  $\mu \in \check{V}_R^*$  (see Figure 3) we get  $\text{Re}(\cosh \mu) < 0$ . Therefore we must choose the sign (-) in (33). So we obtain the first formula of (31).

For  $\mu \in \check{V}_L^*$  (see Figure 3) we obtain  $\text{Re}(\cosh \mu) > 0$ . Hence we get the second formula of (31). The formula (32) is proved similarly.  $\square$

Now we prove the following theorem, which completes the proof of Theorem 3.5.

**Theorem 7.1.** The algebraic equation (16) has no non-trivial solutions in the class  $\mathcal{H}(\mathbb{C}^+) \times \mathcal{H}(\mathbb{C}^+)$ .

*Proof.* Lemma 5.1, the equivalence of the systems (21) and (29) imply that it suffices to prove that the system (29) has no non-trivial solutions in  $\mathcal{H}(\check{V}_1^+) \times \mathcal{H}(\check{V}_2^+)$ .

Let us consider the equation (30) in the right hand side  $\check{V}_R^*$  and in the left hand side  $\check{V}_L^*$  of the  $\check{V}^*$ .

i) We consider (30) in  $\check{V}_R^*$  (see Figure 3). By (31) and (32) the equation (30) is equivalent to  $-2 \cosh \mu \check{v}_1^0(\mu) = 0$ ,  $\mu \in \check{V}_R^*$ , which satisfies if and only if  $\check{v}_1^0(\mu) = 0$ ,  $\mu \in \check{V}_R^*$ . By analyticity  $\check{v}_1^0(\mu)$  in  $\check{V}_{1R}^+$  it follows that  $\check{v}_1^0(\mu) \equiv 0$  in  $\check{V}_{1R}^+$ . Thus by the automorphism  $\check{h}_1$  we conclude  $\check{v}_1^0(\mu) \equiv 0$ ,  $\mu \in \check{V}_1^+$ .

ii) Now we consider (30) in  $\check{V}_L^*$  (see Figure 3). By (31) and (32) the equation (30) is equivalent to  $2i \sinh \mu \check{v}_2^0(\mu) = 0$ ,  $\mu \in \check{V}_L^*$ , which is satisfied only by  $\check{v}_2^0(\mu) = 0$ , for  $\mu \in \check{V}_L^*$ . By analyticity  $\check{v}_2^0(\mu)$  on  $\check{V}_{2L}^+$  it follows that  $\check{v}_2^0(\mu) \equiv 0$  on  $\check{V}_{2L}^+$ . So, by the automorphism  $\check{h}_2$  we conclude  $\check{v}_2^0(\mu) \equiv 0$ ,  $\mu \in \check{V}_2^+$ . The theorem is proved.  $\square$

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## References

- [1] A. Bamberger, P. Joly, J. E. Roberts, Second order absorbing boundary conditions for the wave equation: a solution for the corner problem, *SIAM J. Numer. Anal.*, **27**, No. 2 (1990), 323-352.
- [2] F. Collino, *Analyse Numérique de Modèles de Propagation D'ondes. Application à la Migration et à L'invasion des Données Sismiques*, Ph.D. Thesis, Université Paris IX (1987).
- [3] B. Engquist, E. Majda, Absorbing boundary conditions for the numerical simulation of waves, *Math. Comp.*, **31**, No. 139 (1977), 629-651.
- [4] B. Engquist, E. Majda E, Radiation boundary conditions for acoustic and elastic wave calculations, *Comm. Pure Appl. Math.*, **32** (1979), 313-357.
- [5] D. Givoli, *Numerical Methods for Problems in Infinite Domains*, Elsevier Scientific Publishing Co., Amsterdam (1992).
- [6] P. Joly, *Analyse Numérique et Mathématique de Problèmes Liés a la Propagation D'ondes Acoustiques, Élastiques et Électromagnétiques*, PhD Thesis, Université Paris IX (1987).
- [7] Jürgen Jost, *Compact Riemann Surfaces. An introduction to contemporary Mathematics*, Universitext Springer (1997).

- [8] A.I. Komech, *Linear Partial Differential Equations with Constant Coefficients*, Partial Differential Equations II, Encyclopaedia of Mathematical Science, Volume **31**, Springer-Verlang, Berlin (1994), 127-260.
- [9] A.I. Komech, A. Merzon, P. Zhevandrov, A method of complex characteristics for elliptic problems in angles, and its applications, *Amer. Math. Soc. Transl., Ser. 2*, No. 206, Amer. Math. Soc., Providence, RI (2002), 125-159.
- [10] H.O. Kreiss, Inicial-boundary value problems for hyperbolic systems, *Comm. Pure Appl. Math.* No. 23 (1970), 277-298.
- [11] Zaldivar, Felipe, *Funciones Algebraicas de una Variable Compleja*, Universidad Autónoma Metropolitana, Unidad Iztapalapa, México (1995).

