## Inverse Problems in Transport Theory

#### PLAMEN STEFANOV

ABSTRACT. We study an inverse problem for the transport equation in a bounded domain in  $\mathbb{R}^n$ . Given the incoming flux on the boundary, we measure the outgoing one. The inverse problem is to recover the absorption coefficient  $\sigma_a(x)$  and the collision kernel k(x,v',v) from this data. This paper is a survey of recent results about general k's without assuming that k depends on a reduced number of variables. We present uniqueness results in dimensions  $n \geq 3$  for the time dependent and the stationary problem, and in the time dependent case we study the inverse scattering problem as well. The proofs are constructive and lead to direct procedures for recovering  $\sigma_a$  and k. For n=2 the problem of recovering k is formally determined and we prove uniqueness for small k and a stability estimates.

### 1. Introduction

This paper is a review of the recent progress in the study of inverse problems for the transport equation in  $\mathbb{R}^n$ ,  $n \geq 2$  by the author and M. Choulli [CSt1], [CSt2], [CSt3], [CSt4] and the author and G. Uhlmann [StU]. We are focused here on the case when the collision kernel k introduced below depends on all of its variables x, v', v. There are a lot of works dealing with k's of the form  $k(x, v' \cdot v)$  that is also physically important but we will not discuss those results here.

Define the transport operator T by

$$Tf = -v \cdot \nabla_x f(x, v) - \sigma_a(x, v) f(x, v) + \int_V k(x, v', v) f(x, v') \, dv', \qquad (1-1)$$

where f = f(x, v) represents the density of a particle flow at the point  $x \in X$  moving with velocity  $v \in V$ . Here  $X \subset \mathbb{R}^n$  is a bounded domain with  $C^1$ -boundary,  $V \subset \mathbb{R}^n$  is the velocity space. We assume that V is an open set in Sections 2 and 3, and that  $V = S^{n-1}$  (and dv is replaced by  $dS_v$ ) in Section 4. All of our results in Sections 2 and 3 hold in the case when  $V = S^{n-1}$  with obvious modifications. In fact, the case  $V = S^{n-1}$  leads to some simplifications,

The author was partly supported by NSF Grant DMS-0196440.

for example there is no need to work locally in the open set  $V \setminus \{0\}$  in some cases and the measure  $d\tilde{\xi}$  can be chosen to be  $d\xi$  in Section 3. The coefficient  $\sigma_a(x,v) \geq 0$  above measures the absorption of particles at the point (x,v) due to change of velocity or absorption by the surrounding media. The collision kernel  $k(x,v',v) \geq 0$  is related to the number of particles that change their velocity from v' to v at the point x.

We study inverse problems for both the time-dependent transport equation

$$(\partial_t - T)f = 0 (1-2)$$

and the stationary transport equation

$$Tf = 0. (1-3)$$

Set  $\Gamma_{\pm} = \{(x, v) \in \partial X \times V : \pm n(x) \cdot v > 0\}$ , where n(x) is the outer normal to  $\partial X$  at  $x \in \partial X$ . A typical boundary value problem for (1–2) or (1–3) is associated with boundary conditions

$$f|_{\Gamma_{-}} = f_{-}, \tag{1-4}$$

where  $f_{-}$  is a given function on  $\Gamma_{-}$ , that depends also on t in the first case. On the boundary we measure the outgoing flux  $f_{+}$  generated by a given incoming flux  $f_{-}$ , i.e., we assume that we know that so called *albedo* operator  $\mathcal{A}$  defined by

$$\mathcal{A}: f_{-} \longmapsto f_{+} = f|_{\Gamma_{+}}. \tag{1-5}$$

The inverse problem that we study is this: Can we determine the functions  $\sigma_a$  and k from a knowledge of  $\mathcal{A}$ ? In the time dependent case we study the inverse scattering problem as well—recovery of  $\sigma_a$  and k from the knowledge of the scattering operator S. In order to ensure uniqueness of recovery of  $\sigma_a$ , we assume that  $\sigma_a$  depends on x only, since otherwise there is clearly no uniqueness. Under this assumption and some natural assumptions in the stationary case that guarantee the solvability of the direct problem (1–2), we not only prove uniqueness but give an explicit solution for the inverse problem. Our approach is based on the study of the singularities of the Schwartz kernel of  $\mathcal{A}$  and we show that all information about  $\sigma_a$  and k is contained in those singularities. For large  $\sigma_a$ , those singularities have very small amplitudes and are hard to measure in real applications, so then this approach is of less interest for the practical recovery of  $\sigma_a$  and k. In Section 4 we prove some stability estimates.

#### 2. The Time Dependent Transport Equation

**2.1.** Main results In this section we present inverse problems results about the time dependent transport equation (1–2)

$$\frac{\partial}{\partial t}u(x,v,t) = -v \cdot \nabla_x u(x,v,t) - \sigma_a(x,v)u(x,v,t) + \int_V k(x,v',v)u(x,v',t) \, dv'.$$
(2-1)

We first introduce some terminology and notation. The production rate  $\sigma_p(x, v')$  is defined by

$$\sigma_p(x, v') = \int_V k(x, v', v) \, dv.$$

Following [RS] we say that the pair  $(\sigma_a, k)$  is admissible if

(i)  $0 \le \sigma_a \in L^{\infty}(\mathbb{R}^n \times V)$ ,

(ii) 
$$0 \le k(x, v', \cdot) \in L^1(V)$$
 for a.e.  $(x, v') \in \mathbb{R}^n \times V$  and  $\sigma_p \in L^{\infty}(\mathbb{R}^n \times V)$ .

Throughout this paper we assume that  $\sigma_a$  and k are extended as 0 for  $x \notin X$ . Set  $T_0 = -v \cdot \nabla_x$  with its maximal domain. It is well-known that  $T_0$  is a generator of a strongly continuous group  $U_0(t)f = f(x - tv, v)$  of isometries on  $L^1(X \times V)$  preserving nonnegative functions. Following the notation in [RS], we introduce the operators

$$\begin{split} [A_1f](x,v) &= -\sigma_a(x,v)f(x,v), & T_1 = T_0 + A_1, & D(T_1) = D(T_0), \\ [A_2f](x,v) &= \int_V k(x,v',v)f(x,v')\,dv', & T = T_0 + A_1 + A_2 & D(T) = D(T_0), \end{split}$$

and set  $A = A_1 + A_2$ . Operators  $A_1$  and  $A_2$  are bounded on  $L^1(X \times V)$  and  $T_1$ , T are generators of strongly continuous groups  $U_1(t) = e^{tT_1}$ ,  $U(t) = e^{tT}$ , respectively [RS]. For  $U_1(t)$  we have an explicit formula

$$[U_1(t)f](x,v) = e^{-\int_0^t \sigma_a(x-sv,v) \, ds} f(x-tv,v), \tag{2-2}$$

while for U(t) we have

$$||U(t)|| \le e^{Ct}, \quad C = ||\sigma_p||_{L^{\infty}}.$$
 (2-3)

We work in the Banach space  $L^1(X \times V)$ , so here ||U(t)|| is the operator norm of U(t) in  $\mathcal{L}(L^1(X \times V))$ . It should be mentioned also that U(t) preserves the cone of nonnegative functions for t > 0.

One can define the wave operators associated with T,  $T_0$  by

$$W_{-} = \underset{t \to \infty}{\text{s-lim}} U(t)U_0(-t), \tag{2-4}$$

$$\tilde{W}_{+} = \underset{t \to \infty}{\text{s-}\lim} U_0(-t)U(t). \tag{2-5}$$

If  $W_{-}$  and  $W_{+}$  exist, one can define the scattering operator

$$S = \tilde{W}_{\perp}W_{\perp}$$

as a bounded operator in  $L^1(X \times V)$ . Scattering theory for (1–1) has been developed in [Hej], [Si], [V1] and we refer to these papers (see also [RS]) for sufficient conditions guaranteeing the existence of S. We would like to mention here also [P1], [U], [E], [St], [V2], [CMS]. An abstract approach based on the Limiting Absorption Principle has been proposed in [Mo]. We show below however that S can always be defined as an operator  $S: L^1_c(\mathbb{R}^n \times V \setminus \{0\}) \to L^1_{loc}(\mathbb{R}^n \times V \setminus \{0\})$ .

The inverse scattering problem for (2-1) is the following: Does S determine uniquely  $\sigma_a$ , k? The answer is affirmative if  $\sigma_a$  is independent of v.

THEOREM 2.1 ([CSt1], [CSt2]). Let  $(\sigma_a, k)$ ,  $(\hat{\sigma}_a, \hat{k})$  be two admissible pairs such that  $\sigma_a$ ,  $\hat{\sigma}_a$  do not depend on v and denote by S,  $\hat{S}$  the corresponding scattering operators. Then, if  $S = \hat{S}$ , we have  $\sigma_a = \hat{\sigma}_a$ ,  $k = \hat{k}$ .

The assumption that  $\sigma_a$ ,  $\hat{\sigma}_a$  depend on x only can be relaxed a little by assuming that they depend on x and |v| only. In the general case however, there is no uniqueness. Assume, for example, that k=0. Then

$$Sf = e^{-\int_{-\infty}^{\infty} \sigma_a(x - sv, v) ds} f, \quad (k = 0)$$

$$(2-6)$$

and therefore S determines  $\int_{-\infty}^{\infty} \sigma_a(x-sv,v) ds$  only for any  $v \in V$ . It is easy to see that those integrals do not determine  $\sigma_a$  uniquely. If  $\sigma_a$  is independent of v/|v| however, then S determines the X-ray transform of  $\sigma_a$  and it therefore determines  $\sigma_a$  in this case.

Next we consider the albedo operator  $\mathcal{A}$ . Assume that X is convex and has  $C^1$ -smooth boundary  $\partial X$ . Consider the measure  $d\xi = |n(x) \cdot v| d\mu(x) dv$  on  $\Gamma_{\pm}$ , where  $d\mu(x)$  is the measure on  $\partial X$ . We will solve the problem

$$\begin{cases} (\partial_t - T)u = 0 & \text{in } \mathbb{R} \times X \times V, \\ u|_{\mathbb{R} \times \Gamma_-} = u_-, \\ u|_{t \ll 0} = 0 \end{cases}$$
 (2-7)

for u(t, x, v), where  $u_{-} \in L^{1}_{c}(\mathbb{R}; L^{1}(\Gamma_{-}, d\xi))$ . Problem (2–7) has a unique solution  $u \in C(\mathbb{R}; L^{1}(X \times V))$  and one defines the albedo operator  $\mathcal{A}$  by

$$\mathcal{A}g = u|_{\mathbb{R} \times \Gamma_{+}} \in L^{1}_{loc}(\mathbb{R}; L^{1}(\Gamma_{+}, d\xi)). \tag{2-8}$$

Therefore,  $\mathcal{A}: L^1_{\mathbf{c}}(\mathbb{R}; L^1(\Gamma_-, d\xi)) \to L^1_{\mathrm{loc}}(\mathbb{R}; L^1(\Gamma_+, d\xi))$ . It can be seen that  $\mathcal{A}g$  can be defined more generally for  $g \in L^1(\mathbb{R} \times \Gamma_-, dt \, d\xi)$  with g = 0 for  $t \ll 0$ . We show below that in fact  $\mathcal{A}$  determines S uniquely and conversely, S determines  $\mathcal{A}$  uniquely by means of explicit formulae in case when X is convex. This generalizes earlier results in [AE], [EP], [P2] showing that there is a relationship between S and  $\mathcal{A}$ . To this end, let us define the extension operators  $J_{\pm}$  and the restriction (trace) operators  $R_{\pm}$  as follows. Set

$$\Omega = \{(x, v) \in \mathbb{R}^n \times V \setminus \{0\} : x - tv \in X \text{ for some } t \in \mathbb{R}\},$$
 (2-9)

and define the functions

$$\tau_{\pm}(x,v) = \min\{t \in \mathbb{R} : (x \pm tv, v) \in \Gamma_{\pm}\}, \quad (x,v) \in \Omega. \tag{2-10}$$

Given  $g \in L^1(\mathbb{R} \times \Gamma_{\pm}, dt d\xi)$ , consider the following operators of extension:

$$(J_{\pm}g)(x,v) = \begin{cases} g(\pm \tau_{\pm}(x,v), \ x \pm \tau_{\pm}(x,v)v, \ v) & \text{if } (x,v) \in \Omega, \\ 0 & \text{otherwise.} \end{cases}$$

It is easy to check that  $J_{\pm}: L^1(\mathbb{R} \times \Gamma_{\pm}, dt d\xi) \to L^1(X \times V)$  are isometric. Denote by  $R_{\pm}$  the operator of restriction

$$R_{\pm}f = f|_{\Gamma_+}, \quad f \in C(\mathbb{R}^n \times V).$$

Although  $R_{\pm}$  is not a bounded operator on  $L^1(X \times V)$  (see [Ce1], [Ce2] and Theorem 3.2 below), we can show that  $R_{\pm}U_0(t)f \in L^1(\mathbb{R} \times \Gamma_{\pm}, dt d\xi)$  is well defined for any  $f \in L^1(X \times V)$ . Denote by  $\chi_{\Omega}$  the characteristic function of  $\Omega$ . We establish the following relationships between S and A.

Theorem 2.2 ([CSt1], [CSt2]). Assume that X is convex. Then

- (a)  $Ag = R_+ U_0(t) S J_- g$  for all  $g \in L^1_c(\mathbb{R} \times \Gamma_-, dt d\xi)$ ,
- (b)  $Sf = J_{+}AR_{-}U_{0}(t)f + (1 \chi_{\Omega})f, f \in L_{c}^{1}(\mathbb{R}^{n} \times V \setminus \{0\}),$
- (c) A extends to a bounded operator

$$\mathcal{A}: L^1(\mathbb{R} \times \Gamma_-, dt d\xi) \to L^1(\mathbb{R} \times \Gamma_+, dt d\xi)$$

if and only if S extends to a bounded operator on  $L^1(X \times V)$ .

We decompose  $L^1(\mathbb{R}^n \times V) = L^1(\Omega) \oplus L^1((\mathbb{R}^n \times V) \setminus \Omega)$ . A similar decomposition holds for  $L^1_c(\mathbb{R}^n \times V \setminus \{0\})$ . Then S leaves invariant both spaces, moreover  $S|_{L^1((\mathbb{R}^n \times V) \setminus \Omega)} = \mathrm{Id}$ , so S can be decomposed as a direct sum  $S = S_\Omega \oplus \mathrm{Id}$ . Denote  $\mathcal{R}_{\pm} = R_{\pm}U_0(\cdot) : L^1(\Omega) \to L^1(\mathbb{R} \times \Gamma_{\pm}, dt \, d\xi)$ . Then  $\mathcal{R}_{\pm}$  are isometric and invertible and  $\mathcal{R}_{\pm}^{-1} = \mathcal{J}_{\pm}$  with  $\mathcal{J}_{\pm} : L^1(\mathbb{R} \times \Gamma_{\pm}, dt \, d\xi) \to L^1(\Omega)$ ,  $\mathcal{J}_{\pm}f := J_{\pm}f|_{L^1(\Omega)}$ . Then we can rewrite Theorem 2.2 (a), (b) as follows:

$$\mathcal{A} = \mathcal{R}_{+} S_{\Omega} \mathcal{J}_{-} \qquad \text{on } L^{1}_{c}(\mathbb{R} \times \Gamma_{-}, dt \, d\xi),$$
  
$$S_{\Omega} = \mathcal{J}_{+} \mathcal{A} \mathcal{R}_{-} \qquad \text{on } L^{1}_{c}(\mathbb{R}^{n} \times V \setminus \{0\}).$$

Thus  $\mathcal{A}$  can be obtained from  $S_{\Omega}$  by a conjugation with invertible isometric operators and vice-versa.

An immediate consequence of Theorem 2.2 is that  $\mathcal{A}$  determines  $\sigma_a$  and k uniquely for  $\sigma_a$  independent of v and X convex. In short, in this case the inverse boundary value problem is equivalent to the inverse scattering problem. However, we can prove uniqueness for the inverse boundary value problem for not necessarily convex domains as well independently of Theorems 2.1 and 2.2.

THEOREM 2.3 [CSt1], [CSt2]. Let  $(\sigma_a, k)$ ,  $(\hat{\sigma}_a, \hat{k})$  be two admissible pairs with  $\sigma_a$ ,  $\hat{\sigma}_a$  independent of v. Then, if the albedo operators  $\mathcal{A}$ ,  $\hat{\mathcal{A}}$  on  $\partial X$  coincide, we have  $\sigma_a = \hat{\sigma}_a$ ,  $k = \hat{k}$ .

2.2. Singular decomposition of the fundamental solution and the kernels of S and A. Proof of the main results in section 2.1. The key to proving the uniqueness results above is to study the singularities of the Schwartz kernel of S and respectively A. To this end we will study first the kernel of the

solution operator of the problem (2–7). Given  $(x', v') \in \mathbb{R}^n \times V \setminus \{0\}$ , consider the problem

$$\begin{cases} (\partial_t - T)u = 0 & \text{in } \mathbb{R} \times \mathbb{R}^n \times V, \\ u|_{t \ll 0} = \delta(x - x' - tv)\delta(v - v'), \end{cases}$$
 (2-11)

 $\delta$  being the Dirac delta function. Problem (2–1) has a unique solution

$$u^{\#}(t,x,v,x',v')$$

depending continuously on t with values in  $\mathcal{D}'(\mathbb{R}^n_x \times V_v \times \mathbb{R}^n_{x'} \times V_{v'} \setminus \{0\})$ . We have the following singular expansion of  $u^{\#}$ .

THEOREM 2.4 [CSt1], [CSt2]. Problem (2–11) has the unique solution  $u^{\#} = u_0^{\#} + u_1^{\#} + u_2^{\#}$ , where

$$\begin{split} u_0^{\#} &= e^{-\int_0^{\infty} \sigma_a(x-sv,v) \, ds} \delta(x-x'-tv) \delta(v-v'), \\ u_1^{\#} &= \int_0^{\infty} e^{-\int_0^s \sigma_a(x-\tau v,v) \, d\tau} e^{-\int_0^{\infty} \sigma_a(x-sv-\tau v',v') \, d\tau} \\ &\qquad \qquad \times k(x-sv,v',v) \delta(x-sv-(t-s)v'-x') \, ds, \\ u_2^{\#} &\in C\big(\mathbb{R}; \, L_{\rm loc}^{\infty}(\mathbb{R}^n_{x'} \times V_{v'}; \, L^1(\mathbb{R}^n_x \times V_v))\big). \end{split}$$

The proof of Theorem 2.4 is based on iterating twice Duhamel's formula

$$U(t-r) = U_1(t-r) + \int_r^t U(t-s)A_2U_1(s-r) \, ds$$

and on estimating the remainder term.

To build the scattering theory for the transport equation, we first show that the wave operators  $W_-$ ,  $\tilde{W}_+$  (see (2–4), (2–5)) exist as operators between the spaces

$$W_{-}: L_{\rm c}^{1}(\mathbb{R}^{n} \times V \setminus \{0\}) \longrightarrow L^{1}(X \times V),$$
  
$$\tilde{W}_{+}: L^{1}(X \times V) \longrightarrow L_{\rm loc}^{1}(\mathbb{R}^{n} \times V \setminus \{0\}).$$

Then we define the scattering operator

$$S = \tilde{W}_{+}W_{-} : L_{c}^{1}(\mathbb{R}^{n} \times V \setminus \{0\}) \longrightarrow L_{loc}^{1}(\mathbb{R}^{n} \times V \setminus \{0\}). \tag{2-12}$$

It can be seen that S is well defined on the wider subset  $\{f : \exists t_0 \text{ such that } U_0(t)f = 0 \text{ for } x \in X, t < -t_0\}$  (the incoming space).

The distribution  $u^{\#}(t, x, v, x', v')$  is the Schwartz kernel of  $U(t)W_{-}$ . Let S(x, v, x', v') be the Schwartz kernel of the scattering operator S. Then

$$S(x, v, x', v') = \lim_{t \to \infty} u^{\#}(t, x+tv, v, x', v').$$

This limit exists trivially, in fact for any  $K \in \mathbb{R}^n \times V \setminus \{0\}$ , the distribution  $u^{\#}(t, x+tv, v, x', v')|_{K}$  is independent of t for  $t \geq t_0(K)$ . On the other hand, as mentioned in the Introduction, it is not trivial to show that under some

condition, S is a kernel of a bounded operator in  $L^1(\mathbb{R}^n \times V)$ . One can also prove the following integral representation of the scattering kernel:

$$S(x, v, x', v') = e^{-\int_{-\infty}^{\infty} \sigma_a(x - \tau v, v) d\tau} \delta(x - x') \delta(v - v')$$

$$+ \int_{-\infty}^{\infty} e^{-\int_s^{\infty} \sigma_a(x + \tau v, v) d\tau} (A_2 u^{\#})(s, x + sv, v, x', v') ds. \quad (2-13)$$

This formula is an analogue of the representation of the scattering amplitude (in our setting, that would be the second term in the right-hand side above) for the Schrödinger equation.

Now, combining Theorem 2.4 and the representation above, we get the following express for the kernel S of the scattering operator S.

THEOREM 2.5 [CSt1], [CSt2]. We have  $S = S_0 + S_1 + S_2$ , where the Schwartz kernels  $S_j(x, v, x', v')$  of the operators  $S_j$ , j = 0, 1, 2 satisfy

$$S_0 = e^{-\int_{-\infty}^{\infty} \sigma_a(x-\tau v,v)d\tau} \delta(x-x')\delta(v-v'),$$

$$S_1 = \int_{-\infty}^{\infty} e^{-\int_s^{\infty} \sigma_a(x+\tau v,v)d\tau} e^{-\int_0^{\infty} \sigma_a(x+sv-\tau v',v')d\tau} k(x+sv,v',v)\delta(x-x'+s(v-v')) ds,$$

$$S_2 \in L^{\infty}_{loc}(\mathbb{R}^n_{x'} \times V_{v'} \setminus \{0\}; L^1_{loc}(\mathbb{R}^n_x \times V_v \setminus \{0\})).$$

We are ready now to complete the proof of Theorem 2.1. The idea of the proof is the following. Suppose we are given the scattering operator S corresponding to a unknown admissible pair  $(\sigma_a, k)$ . Then we know the kernel  $\mathcal{S} = \mathcal{S}_0 + \mathcal{S}_1 + \mathcal{S}_2$ . It follows from Theorem 2.5 that  $S_0$  is a singular distribution supported on the hyperplane x = x', v = v' of dimension 2n,  $S_1$  is supported on a (3n + 1)dimensional surface (for  $v \neq v'$ ), while  $S_2$  is a function. Therefore,  $S_j$ , j = 0, 1, 2have different degrees of singularities and given  $S = S_0 + S_1 + S_2$ , one can always recover  $S_0$  and  $S_1$ . From  $S_0$  one can recover the X-ray transform of  $\sigma_a$ and therefore  $\sigma_a$  itself, provided that  $\sigma_a$  is independent of v. Next, suppose for simplicity that  $\sigma_a$  and k are continuous. Then for fixed x, v, v' with  $v \neq v', \mathcal{S}_1$ is a delta-function supported on the line x' = x + s(v - v'),  $s \in \mathbf{R}$  with density a multiple of k(x+sv,v',v). Therefore, one can recover that density for each s and in particular setting s=0, we get k(x,v',v). Moreover, based on this, we can write explicit formulae that extract  $\sigma_a$  and k from S by allowing S to act on a sequence of suitably chosen test functions that concentrate on one of the singular varieties described above, see [CSt2].

We will skip the proof of Theorem 2.2. We merely recall that it proves uniqueness for the inverse boundary value problem for convex X as a direct consequence of Theorem 2.5.

Now assume that X is not necessarily convex. We can still prove uniqueness for the inverse boundary value problem by showing that  $\mathcal{A}$  determines uniquely  $u^{\#}$  for x outside X by following arguments in [SyU], and then using (2–13) and Theorem 2.5. In order to give a constructive (in fact, explicit) reconstruction, we

study next the Schwartz kernel of the operator  $\mathcal{A}$  in the spirit of Theorem 2.5. A priori, this kernel is a distribution in  $\mathcal{D}'(\mathbb{R} \times \Gamma_+ \times \mathbb{R} \times \Gamma_-)$ . Denote by  $\delta_1$  the Dirac delta function on  $\mathbb{R}^1$  and by  $\delta_y(x)$  the Delta function on  $\partial X$  defined by  $\int_{\partial X} \delta_y \varphi \, d\mu(x) = \varphi(y)$ . Set

$$\theta(x,y) = \begin{cases} 1 & \text{if } px + (1-p)y \in X \text{ for each } p \in [0,1], \\ 0 & \text{otherwise.} \end{cases}$$

Theorem 2.6 [CSt1], [CSt2]. The Schwartz kernel of A has the form

$$\alpha(t-t',x,v,x',v'),$$

that is, formally,

$$(\mathcal{A}g)(t,x,v) = \int_{\mathbb{R}\times\Gamma_{-}} \alpha(t-t',x,v,x',v')g(t',x',v')\,d\xi(x',v')$$

with  $\alpha = \alpha_0 + \alpha_1 + \alpha_2$ , where the  $\alpha_j(\tau, x, v, x', v')$  (with  $(x, v) \in \Gamma_+$ ,  $(x', v') \in \Gamma_-$ ) satisfy

$$\alpha_{0} = |n(x') \cdot v'|^{-1} e^{-\int_{0}^{\tau_{-}(x,v)} \sigma_{a}(x-sv,v) ds} \delta_{\{x-\tau_{-}(x,v)v\}}(x') \delta(v-v') \delta_{1}(\tau-\tau_{-}(x,v)),$$

$$\alpha_{1} = |n(x') \cdot v'|^{-1} \int e^{-\int_{0}^{s} \sigma_{a}(x-pv,v) dp} e^{-\int_{0}^{\tau_{-}(x-sv,v')} \sigma_{a}(x-sv-pv',v') dp} \times$$

$$\delta_{1}(\tau-s-\tau_{-}(x-sv,v')) k(x-sv,v',v) \delta_{\{x-sv-\tau_{-}(x-sv,v')v'\}}(x') \theta(x-sv,x) ds,$$

$$\alpha_{2} \in L^{\infty} \left(\Gamma_{-}; L^{1}_{loc}(\mathbb{R}_{\tau}; L^{1}(\Gamma_{+},d\xi))\right).$$

The proof of Theorem 2.3 now follows from the theorem above by analyzing the singularities of  $\alpha$  as above. In this case, we can also wite explicit formulae for  $\sigma_a$  and k as certain limits of the distribution  $\alpha$  acting on a sequence of test functions concentrating near the singularities of  $\alpha_0$  and  $\alpha_1$ , respectively; see [CSt2].

#### 3. The Stationary Transport Equation

We turn our attention now to the boundary value problem (1–3), (1–4) for the stationary transport equation

$$\begin{cases} -v \cdot \nabla_x f(x, v) - \sigma_a(x, v) f(x, v) + \int_V k(x, v', v) f(x, v') \, dv' = 0 & \text{in } X \times V, \\ f|_{\Gamma_-} = f_-. & (3-1) \end{cases}$$

Here X does not need to be convex. We impose some conditions in order to ensure the unique solvability of the direct problem (3–1). Recall the definition (2–10) of  $\tau_{\pm}$ . Set  $\tau = \tau_{-} + \tau_{+}$ . We will consider two cases. First we assume that

$$\|\tau\sigma_a\|_{L^\infty} < \infty, \qquad \|\tau\sigma_p\|_{L^\infty} < 1. \tag{3-2}$$

This condition in particular guarantees that the dynamics U(t) corresponding to the time-dependent transport equation  $(\partial_t - T)u = 0$  is subcritical [RS], that is, the  $L^1$ -norm of the solution is uniformly bounded for t > 0; compare with (2–3).

Note that (3–2) holds if in particular  $\||v|^{-1}\sigma_a\|_{L^{\infty}} < \infty$ , diam  $X\||v|^{-1}\sigma_p\|_{L^{\infty}} < 1$ . The second situation we will consider is when [DL]

$$0 \le \nu \le \sigma_a(x, v) - \sigma_p(x, v)$$
 for a.e.  $(x, v) \in X \times V$  (3-3)

with some  $\nu > 0$ . This condition means that the absorption rate is greater than the production rate. This also implies that the corresponding dynamics is subcritical.

The main result in this section is this:

THEOREM 3.1 [CSt3], [CSt4]. Assume that  $(\sigma_a, k)$ ,  $(\hat{\sigma}_a, \hat{k})$  are two admissible pairs with  $\sigma_a = \sigma_a(x, |v|)$ ,  $\hat{\sigma}_a = \hat{\sigma}_a(x, |v|)$  and assume that they satisfy either (3–2) or (3–3). Assume that the corresponding albedo operators  $\mathcal{A}$  and  $\hat{\mathcal{A}}$  coincide. Then

- (a) if  $n \geq 3$ , then  $\sigma_a = \hat{\sigma}_a$ ,  $k = \hat{k}$ ;
- (b) if n = 2, then  $\sigma_a = \hat{\sigma}_a$ .

In the stationary case we have less data than in the time dependent one studied in Section 2 because the time variable is not present. The inverse boundary value problem is overdetermined in dimension  $n \geq 3$  because the kernel of  $\mathcal{A}$  depends on 4n-2 variables while k is a function of 3n variables,  $\sigma_a(x,v/|v|)$  is a function of 2n-1 variables. In the 2D case however, the inverse problem for recovery of k is formally determined (but still overdetermined for the recovery of  $\sigma_a$ ) and the theorem above does not provide uniqueness in this case. In Section 4 we formulate a uniqueness result in the 2D case for small k.

To prove Theorem 3.1, we again study the Schwartz kernel of the albedo operator  $\mathcal{A}$ . Note that in the stationary case,  $\mathcal{A}$  acts on functions on  $\Gamma_{-}$  and maps them to functions on  $\Gamma_{+}$  (see next subsection for details). It turns out that, similarly to the time dependent case, the kernel of  $\mathcal{A}$  has a singular decomposition  $\alpha = \alpha_0 + \alpha_1 + \alpha_2$ , where  $\alpha_0$  and  $\alpha_1$  are delta functions for  $n \geq 3$  supported on varieties of different dimensions with densities that identify respectively  $\sigma_a$  and k. The third term  $\alpha_2$  is a locally  $L^1$  function and can be distinguished form  $\alpha_0$  and  $\alpha_1$ . If n = 2,  $\alpha_0$  is still a delta function, but  $\alpha_1$  is a locally  $L^1$  function and cannot be distinguished from  $\alpha_2$  and this explains why this method does not work then. The 2D case is considered separately in next section.

3.1. Some estimates in the stationary case Similarly to the measure defined in Section 2.1, we define a measure on  $\Gamma_{\pm}$  by:

$$d\tilde{\xi} = \min\{\tau(x, v), \lambda\} |n(x) \cdot v| d\mu(x) dv,$$

where  $\lambda > 0$  is an arbitrary constant. Using the trace Theorem 3.2 below, we can show that  $\mathcal{A}: L^1(\Gamma_-, d\xi) \to L^1(\Gamma_+, d\xi)$  is a bounded operator if (3–2) holds and  $\mathcal{A}: L^1(\Gamma_-, d\tilde{\xi}) \to L^1(\Gamma_+, d\tilde{\xi})$  is bounded if (3–3) holds; see also [DL]. Note that if a neighborhood of the origin is not included in V or in particular, if  $V = S^{n-1}$ , then we can choose  $d\tilde{\xi} = \tau d\xi$ .

Theorem 3.2 (Trace Theorem). (a) [CSt3]

$$||f|_{\Gamma_{\pm}}||_{L^1(\Gamma_+,d\xi)} \le ||T_0f|| + ||\tau^{-1}f||.$$

(b) [Ce1], [Ce2]

$$||f|_{\Gamma_{\pm}}||_{L^{1}(\Gamma_{+},d\tilde{\xi})} \le \lambda ||T_{0}f|| + ||f||.$$

Note that (b) follows from (a) by setting  $f = \min\{\lambda, \tau\}g$ .

We introduce the spaces W,  $\tilde{W}$  via the norms

$$||f||_{\mathcal{W}} = ||T_0 f|| + ||\tau^{-1} f||, \qquad ||f||_{\tilde{\mathcal{W}}} = ||T_0 f|| + ||f||.$$

Then Theorem 3.2 says that taking the trace  $f \mapsto f|_{\Gamma_{\pm}}$  is a continuous operator from W into  $L^1(\Gamma_{\pm}, d\xi)$  and similarly from  $\tilde{W}$  into  $L^1(\Gamma_{\pm}, d\tilde{\xi})$ .

Given  $f_- \in L^1(\Gamma_-, d\xi)$ , define  $Jf_-$  as the following prolongation of  $f_-$  inside  $X \times V$ :

$$Jf_{-} = e^{-\int_{0}^{\tau_{-}(x,v)} \sigma_{a}(x-sv,v) \, ds} f_{-}(x-\tau_{-}(x,v)v,v), \qquad (x,v) \in X \times V.$$

 $Jf_-$  is defined so that  $T_1Jf_-=0$ ,  $Jf_-|_{\Gamma_-}=f_-$ , therefore J is the solution operator of the problem  $T_1f=0$ ,  $f|_{\Gamma_-}=f_-$ .

Proposition 3.1. (a) Assume that  $\|\tau\sigma_a\|_{L^{\infty}} < \infty$ . Then

$$||Jf_-||_{\mathcal{W}} \le C||f_-||_{L^1(\Gamma_-,d\xi)},$$

with  $C = 1 + \|\tau \sigma_a\|_{L^{\infty}}$ . If  $\sigma_a = 0$ , then we have equality above (and C = 1). (b) Assume (3-3). For any  $f_- \in L^1(\Gamma_-, d\tilde{\xi})$  we have

$$||Jf_-||_{\tilde{\mathcal{W}}} \le C||f_-||_{L^1(\Gamma_-,d\tilde{\mathcal{E}})},$$

where 
$$C = (1 + \|\sigma_a\|_{L^{\infty}}) \max\{1, (\nu\lambda)^{-1}\}.$$

We reduce the boundary value problem (3–1) to an integral equation using standard arguments to get

$$(I+K)f = Jf_{-}, (3-4)$$

where I stands for the identity and K is the integral operator

$$Kf = -\int_0^{\tau_-(x,v)} e^{-\int_0^t \sigma_a(x-sv,v) \, ds} (A_2 f)(x-tv,v) \, dt.$$
 (3-5)

Formally,  $K = T_1^{-1}A_2$ , and for  $T_1^{-1}$  we have

$$T_1^{-1}f = -\int_0^{\tau_-(x,v)} e^{-\int_0^t \sigma_a(x-sv,v) \, ds} f(x-tv,v) \, dt.$$

Proposition 3.2. Assume (3–2). Then:

(a)  $\tau^{-1}T_1^{-1}$ ,  $\tau^{-1}T^{-1}$  and  $A_2\tau$  are bounded operators in  $L^1(X \times V)$  and therefore  $K = T_1^{-1}A_2$  is a bounded operator in  $L^1(X \times V; \tau^{-1} \, dx \, dv)$ . Moreover, the operator norm of K is not greater than  $\|\tau\sigma_p\|_{L^\infty} < 1$  and therefore  $(I+K)^{-1}$  exists in this space.

- (b) The integral equation (3–4) and therefore the boundary value problem (3–1) are uniquely solvable for any  $f_- \in L^1(\Gamma_-, d\xi)$  and then  $f \in \mathcal{W}$ .
- (c) The albedo operator  $\mathcal{A}$  is a bounded map  $\mathcal{A}: L^1(\Gamma_-, d\xi) \to L^1(\Gamma_+, d\xi)$ .

Proposition 3.3. Assume (3–3). Then:

- (a) K,  $T_1^{-1}$  and  $T^{-1}$  are bounded operators in  $L^1(X \times V)$  and  $K = T_1^{-1}A_2$ . Further, I + K is invertible and  $(I + K)^{-1} = I - T^{-1}A_2$ .
- (b) The integral equation (3–4) and therefore the boundary value problem (3–1) are uniquely solvable for any  $f_- \in L^1(\Gamma_-, d\tilde{\xi})$  and then  $f \in \tilde{\mathcal{W}}$ .
- (c)  $\mathcal{A}$  is a bounded map  $\mathcal{A}: L^1(\Gamma_-, d\tilde{\xi}) \to L^1(\Gamma_+, d\tilde{\xi})$ .

# 3.2. The fundamental solution of the stationary transport equation. We solve (3–1) with $f_- = \phi_-$ , where

$$\phi_{-} = \frac{1}{|n(x') \cdot v'|} \delta_{\{x'\}}(x) \delta(v - v'),$$

where  $(x',v') \in \Gamma_-$  are regarded as parameters, n(x') is the outer normal, and  $\delta_{\{x'\}}$  is a distribution on  $\partial X$  defined by  $(\delta_{\{x'\}},\varphi) = \int \delta_{\{x'\}}(x)\varphi(x)d\mu(x) = \varphi(x')$ . On the other hand, we will denote by  $\delta$  the ordinary Dirac delta function in  $\mathbb{R}^n$ . The integral above is to be considered in distribution sense. Let us denote by  $\phi(x,v,x',v')$  the solution (in distribution sense) of

$$\begin{cases}
T\phi = 0 & \text{in } X \times V, \\
\phi|_{\Gamma_{-}} = \phi_{-}.
\end{cases}$$
(3-6)

To solve (3–6), we write

$$\varphi = J\varphi_{-} - KJ\varphi_{-} + (I+K)^{-1}K^{2}J\varphi_{-}$$

and analyze each term. To treat the third term, we observe that

$$(I+K)^{-1}K^2J\varphi_- = T^{-1}A_2KJ\varphi_-.$$

We are led to:

THEOREM 3.3. Assume that  $(\sigma_a, k)$  is admissible and either (3–2) or (3–3) holds. Then the solution  $\phi(x, v, x', v')$  of (3–6) can be written as  $\phi = \phi_0 + \phi_1 + \phi_2$ , where

$$\begin{split} \phi_0 &= \int_0^{\tau_+(x',v')} e^{-\int_0^{\tau_-(x,v)} \sigma_a(x-pv,v) \, dp} \delta(x-x'-tv) \delta(v-v') \, dt, \\ \phi_1 &= \int_0^{\tau_-(x,v)} \int_0^{\tau_+(x',v')} e^{-\int_0^s \sigma_a(x-pv,v) \, dp} e^{-\int_0^{\tau_-(x-sv,v')} \sigma_a(x-sv-pv',v') \, dp} \\ &\qquad \qquad \times k(x-sv,v',v) \delta(x-x'-sv-tv') \, dt \, ds, \end{split}$$

and

$$\phi_2 \in L^{\infty}(\Gamma_-; \mathcal{W}) \qquad if (3-2) \ holds,$$
$$\min\{\tau, \lambda\}^{-1}\phi_2 \in L^{\infty}(\Gamma_-; \tilde{\mathcal{W}}) \qquad if (3-3) \ holds.$$

The so constructed solution  $\phi(x, v, x', v')$  is the distribution kernel of the solution operator  $f_- \mapsto f$  of (3–1). In order to find the distribution kernel  $\alpha(x, v, x', v')$   $((x, v) \in \Gamma_+, (x', v') \in \Gamma_-)$  of the albedo operator  $\mathcal{A}$ , it is enough to set

$$\alpha(x, v, x', v') := \phi(x, v, x', v')|_{(x,v) \in \Gamma_+}, \qquad (x', v') \in \Gamma_-.$$

Then, in the sense of distributions,

$$(\mathcal{A}f_{-})(x,v) = \int_{\Gamma} \alpha(x,\theta,x',\theta') f_{-}(x',\theta') d\xi(x',\theta') \quad \text{for all } f_{-} \in C_{0}^{\infty}(\Gamma_{+}).$$

Theorem 3.3 yields the following.

THEOREM 3.4. Assume that  $(\sigma_a, k)$  is an admissible pair and that either (3–2) or (3–3) holds. Then the distribution kernel  $\alpha(x, v, x', v')$  of  $\mathcal{A}$  satisfies  $\alpha = \alpha_0 + \alpha_1 + \alpha_2$  with

$$\alpha_0 = \frac{1}{n(x) \cdot v} e^{-\int_0^{\tau_-(x,v)} \sigma_a(x-pv,v) dp} \delta_{\{x'+\tau_+(x',v')v'\}}(x) \delta(v-v'),$$

$$\alpha_1 = \frac{1}{n(x) \cdot v} \int_0^{\tau_+(x',v')} e^{-\int_0^{\tau_+(x'+tv',v)} \sigma_a(x-pv,v) dp} e^{-\int_0^t \sigma_a(x+pv',v') dp} \times k(x+tv',v',v) \delta_{\{x'+tv'+\tau_+(x'+tv',v)\}}(x) dt,$$

and

$$\alpha_2 \in L^{\infty}(\Gamma_-; L^1(\Gamma_+, d\xi)) \qquad if (3-2) \ holds,$$
$$\min\{\tau(x', v'), \lambda\}^{-1}\alpha_2 \in L^{\infty}(\Gamma_-; L^1(\Gamma_+, d\tilde{\xi})) \qquad if (3-3) \ holds.$$

Theorem 3.4 implies the following way for proving Theorem 3.1. Assume that we are given the albedo operator  $\mathcal{A}$ , corresponding to some admissible pair  $(\sigma_a, k)$ , satisfying either (3–2) or (3–3). Then we also know the distribution  $\alpha(x, v, x', v')$ . By Theorem 3.4,  $\alpha = \alpha_0 + \alpha_1 + \alpha_2$ . Here  $\alpha_0$  is a delta-type distribution supported on a (2n-1)-dimensional variety in  $\Gamma_+ \times \Gamma_-$ . Next,  $\alpha_1$  is also a delta-type distribution (provided that  $n \geq 3$ ) supported on a 3n-dimensional variety in  $\Gamma_+ \times \Gamma_-$ , while  $\alpha_2$  is a (locally  $L^1$ ) function on the (4n-2)-dimensional  $\Gamma_+ \times \Gamma_-$ . Notice that if n=2, then  $\alpha_1$  is a function as well. Therefore, if  $n\geq 3$ , one can distinguish between  $\alpha_0 + \alpha_1$  and  $\alpha_2$ . Moreover, since  $\alpha_0$  and  $\alpha_1$  have different degrees of singularities, one can recover  $\alpha_0$  and  $\alpha_1$ . Now, if  $\sigma_a = \sigma_a(x, |v|)$ , then  $\alpha_0$  determines the X-ray transform  $\int \sigma_a(x+s\omega,|v|) ds$  of  $\sigma_a$  for all x, |v| and  $\omega$  in an open subset of  $S^{n-1}$  (for all  $\omega \in S^{n-1}$  if V is spherically symmetric). This determines uniquely  $\sigma_a$  (see e.g. [H]). Next, once we know  $\sigma_a$ , from  $\alpha_1$ we can recover k. If n=2, then we can recover  $\alpha_0$  and therefore  $\sigma_a$ , but we cannot (at least using those arguments) distinguish between  $\alpha_1$  and  $\alpha_2$  which are both functions and therefore this approach does not work for reconstructing k in two dimensions. Based on those arguments, we can write explicit formulas for recovering the X-ray transform of  $\sigma_a$  and for recovering k as limits of the action of  $\alpha$  on certain sequences of test functions with supports shrinking to the singularities of  $\alpha_0$  and  $\alpha_1$ , respectively; see [CSt4].

#### 4. The Two-Dimensional Stationary Transport Equation

In this section we study the inverse problem for the stationary transport equation (1–3) in the 2D case. As explained in Section 2, the inverse problem of recovering k from the albedo operator  $\mathcal{A}$  is formally determined in this case. We prove below that there exists unique solution provided that k is small enough in the  $L^{\infty}$  norm and we also derive a stability estimate. The results of this section are based on joint work by the author and G. Uhlmann [StU].

Let  $X \subset \mathbb{R}^2$  be an open convex set with smooth boundary and let us write the 2D stationary transport equation (1–3) in the form

$$-v(\theta) \cdot \nabla_x f - \sigma_a(x,\theta)f + \int_{S^1} k(x,\theta',\theta)f(x,\theta') d\theta' = 0, \qquad (4-1)$$

where  $x \in X$ ,  $\theta, \theta' \in S^1$  and

$$v(\theta) = (\cos \theta, \sin \theta).$$

We assume that  $\sigma_a$  and k are in  $L^{\infty}$ . For simplicity, we consider the case  $V = S^1$  here.

In this paper we work with k sufficiently small and under this assumption the direct problem (3–1) is always uniquely solvable with  $f_- \in L^{\infty}(X \times S^1)$  and  $\mathcal{A}f_- \in L^{\infty}(\Gamma_+)$ . For the purpose of the inverse problem however, it is enough to think of  $\mathcal{A}$  as an operator mapping  $C_0^{\infty}(\Gamma_-)$  to  $L^{\infty}(\Gamma_+)$ .

Uniqueness and stability in the 2D case are proved for small k in the case when  $k = k(x, \cos(\theta - \theta'))$  in [R1], and in the case  $k = k(\theta, \theta')$  uniqueness for small k is proved in [T]. Note that in those cases the inverse problem is still formally overdetermined. In this section we prove a uniqueness results for general (but still small)  $k(x, \theta', \theta)$ . Our first result addresses the uniqueness of this inverse problem.

Theorem 4.1 [StU]. Define the class

$$\mathcal{U}_{\Sigma,\varepsilon} = \left\{ (\sigma_a(x), k(x, \theta', \theta)) : \|\sigma_a\|_{L^{\infty}} \le \Sigma, \|k\|_{L^{\infty}} \le \varepsilon \right\}. \tag{4-2}$$

For any  $\Sigma > 0$  there exists  $\varepsilon > 0$  such that a pair  $(\sigma_a, k) \in \mathcal{U}_{\Sigma, \varepsilon}$  is uniquely determined by its albedo operator  $\mathcal{A}$  in the class  $\mathcal{U}_{\Sigma, \varepsilon}$ .

To prove Theorem 4.1, we study again singularities of the distribution kernel  $\alpha(x,\theta,x',\theta')$  of  $\mathcal{A}$  as in Section 3. In the two dimensional case however  $\alpha=\alpha_0+\alpha_1+\alpha_2$  with  $\alpha_0$ ,  $\alpha_1$ ,  $\alpha_2$  as in Theorem 3.4, but  $\alpha_1$  is not a delta type distribution anymore, instead it is an  $L^1$  function and cannot be distinguished from  $\alpha_2$  as in Section 3. We denote  $b=\alpha_1+\alpha_2$ . The term  $\alpha_1$  does have a singularity (integrable) at  $\theta=\pm\theta'$  and it is  $L^\infty_{\rm loc}$  outside this set. Similarly,  $\alpha_2$  has a weaker, logarithmic singularity, as shown below. We will prove below that  $\sin(\theta-\theta')b\in L^\infty$ . We can therefore write

$$\alpha = \frac{\delta(\theta - \theta')\delta_{\{x' + \tau_+(x', \theta')v(\theta')\}}(x)}{n(x) \cdot v(\theta)} e^{-a(x', \theta')} + b(x, \theta, x', \theta'), \tag{4-3}$$

where

$$a(x', \theta') = \int_0^{\tau_+(x', \theta')} \sigma_a(x' + tv(\theta'), \theta') dt, \qquad \sin(\theta - \theta')b \in L^{\infty}. \tag{4-4}$$

In particular, knowing  $\mathcal{A}$ , we can uniquely determine a and b. Let we have two pairs  $(\sigma_a, k)$  and  $(\tilde{\sigma}_a, \tilde{k})$  with albedo operators  $\mathcal{A}$  and  $\tilde{\mathcal{A}}$ , respectively. Set

$$\delta_1 = \|a - \tilde{a}\|_{H^1(\Gamma_-)}, \qquad \delta_2 = \|(b - \tilde{b})\sin(\theta - \theta')\|_{L^{\infty}(\Gamma_- \times \Gamma_+)}.$$
 (4-5)

Our second result is the following stability estimate.

Theorem 4.2 [StU]. Let

$$\mathcal{V}_{\Sigma,\varepsilon}^{s} = \left\{ (\sigma_{a}(x), k(x, \theta', \theta)) \in H^{s}(X) \times C(X \times S^{1} \times S^{1}) : \|\sigma_{a}\|_{H^{s}} \leq \Sigma, \|k\|_{L^{\infty}} \leq \varepsilon \right\}. \quad (4-6)$$

Then, for any s > 1,  $\Sigma > 0$ , there exists  $\varepsilon > 0$  such that for any  $(\sigma_a, k) \in \mathcal{V}_{\Sigma, \varepsilon}^s$  and  $(\tilde{\sigma}_a, \tilde{k}) \in \mathcal{V}_{\Sigma, \varepsilon}^s$  and  $0 < \mu < 1 - 1/s$ , there exists C > 0 such that

$$\|\sigma_a - \tilde{\sigma}_a\|_{L^{\infty}} \le C\delta_1^{1-1/s-\mu},$$
  
 $\|k - \tilde{k}\|_{L^{\infty}} \le C(\delta_1^{1-1/s-\mu} + \delta_2).$ 

REMARK. It follows from (4–15) and (4–16) below that we can choose  $\varepsilon = C(d)e^{-2d\Sigma}$ , where  $d = \operatorname{diam} D$ .

IDEA OF THE PROOF OF THEOREM 4.1. Recall that  $\sigma_a$  and k are  $L^{\infty}$  in all variables. It is convenient to think later that  $\sigma_a$  and k are extended as 0 for  $x \notin X$ .

First we reduce the boundary value problem

$$\begin{cases}
Tf = 0 & \text{in } X \times S^1, \\
f|_{\Gamma_-} = f_- \in C_0^{\infty}(\Gamma_+)
\end{cases}$$
(4-7)

to the integral equation (3-4). Then f is given by

$$f = (I+K)^{-1}Jf_{-}, (4-8)$$

provided that I+K is invertible in a suitable space. The definition (3–5) of K implies immediately that

$$||Kf||_{L^{\infty}(X\times S^{1})} \le C||f||_{L^{\infty}(X\times S^{1})},$$

where  $C = \operatorname{diam} X \|k\|_{L^{\infty}}$ . Therefore, if  $(\sigma_a, k) \in \mathcal{U}_{\Sigma,\varepsilon}$  and  $\varepsilon < 1/\operatorname{diam} X$ , then I + K is invertible in  $L^{\infty}(X \times S^1)$ , and then the solution f to (4–7) is given by (4–8). By using Neumann series, it is not hard to see that the trace  $(I+K)^{-1}f|_{\Gamma_+}$  is well defined in  $L^{\infty}(\Gamma_+)$  for any  $f \in L^{\infty}(X \times S^1)$ . This proves in particular that  $\mathcal{A}$  maps  $C_0^{\infty}(\Gamma_-)$  into  $L^{\infty}(\Gamma_+)$  under the smallness assumption on k above. The same arguments also show that  $\mathcal{A}f_-$  can be defined for any  $f_- \in L^{\infty}(\Gamma_-)$  as well but we will not need this since we work with the distribution kernel of  $\mathcal{A}$ .

Define the fundamental solution  $\phi(x, \theta, x', \theta')$  of the boundary value problem (4–7) as in Section 3. For  $(x', v') \in \Gamma_-$ , let  $\phi(x, v, x', v')$  solve

$$\begin{cases}
T\phi = 0 & \text{in } X \times S^1, \\
\phi|_{\Gamma_-} = |n(x') \cdot v(\theta')|^{-1} \delta_{x'}(x) \delta(\theta - \theta').
\end{cases}$$
(4-9)

As before, the albedo operator  $\mathcal{A}$  has distribution kernel

$$\alpha(x, \theta, x', \theta') = \phi(x, \theta, x', \theta')|_{(x, \theta) \in \Gamma_+},$$

with  $(x', \theta') \in \Gamma_-, (x, \theta) \in \Gamma_+.$ 

As in Section 3, we construct a singular expansion  $\phi = \phi_0 + \phi_1 + \phi_2$  as follows. Let

$$E(x, \theta, t) = e^{\mp \int_0^t \sigma_a(x+sv(\theta), \theta) ds}, \quad \pm t > 0$$

be the total absorption along the path  $[x, x+tv(\theta)]$ . Then

$$\phi_0 = J\phi_-, \qquad \phi_1 = KJ\phi_-, \qquad \phi_2 = (I+K)^{-1}K^2J\phi_-,$$
 (4-10)

and  $\phi_{-} = |n(x') \cdot v(\theta')|^{-1} \delta_{x'}(x) \delta(\theta - \theta')$  as in (4-9). Next,

$$\phi_0 = E(x, \theta, -\infty)\delta(\theta - \theta') \int_0^{\tau_+(x', v(\theta'))} \delta(x - x' - tv(\theta')) dt$$

and

$$\phi_1 = \chi E(y, \theta', -\infty) \frac{k(y, \theta', \theta)}{|\sin(\theta - \theta')|} E(y, \theta, \infty), \tag{4-11}$$

where  $y = y(x', \theta', x, v)$  is the point of intersection of the rays  $(0, \infty) \ni s \mapsto x' + sv(\theta')$  and  $(-\infty, 0) \ni t \mapsto x + tv(\theta)$  and  $\chi = \chi(x, \theta, x', \theta')$  equals 1, if those two rays intersect in  $\bar{X}$ , otherwise  $\chi = 0$ . Recall that X is convex.

To estimate  $\phi_2$ , we need a lemma.

LEMMA 4.1. Let  $(\sigma_a, k)$  and  $(\sigma_a, \tilde{k})$  be in  $L^{\infty}$ . Let  $A_2$ , K and  $\tilde{A}_2$ ,  $\tilde{K}$  be related to k and  $\tilde{k}$  (not necessarily nonnegative), respectively. Then there exists C > 0 depending on diam X only such that

$$|(\tilde{K}K\phi_0)(x,\theta,x',\theta')| \le C||\tilde{k}||_{L^{\infty}}||k||_{L^{\infty}}\left(1 + \log\frac{1}{\sin|\theta - \theta'|}\right)$$

almost everywhere on  $X \times S^1 \times \Gamma_-$ , and also almost everywhere on  $\Gamma_+ \times \Gamma_-$ .

The proof of this lemma is based on the estimate

$$\left| (\tilde{A}_2 K \phi_0)(x, \theta, x', \theta') \right| \le \int_{y \in l(x', \theta')} \frac{\left| \tilde{k}(x, \arg(x-y), \theta) k(y, \theta', \arg(x-y)) \right|}{|x-y|} dl(y),$$

where  $l(x', \theta')$  is the line through x' parallel to  $\theta'$  and dl is the Euclidean measure on it. Using the following elementary estimate

$$\int_{-A}^{A} \frac{ds}{\sqrt{\nu^2 + s^2}} \le 2\left(1 + \log\frac{A}{\nu}\right), \qquad 0 < \nu \le A,$$

we easily complete the proof of the lemma.

For  $\phi_2$  we therefore have  $\phi_2 = (I + K)^{-1}\phi_2^{\#}$  with  $\phi_2^{\#} = K^2\phi_0$  and by Lemma 4.1,

$$0 \le \phi_2^{\#}(x', \theta', x, \theta) \le C \|k\|_{L^{\infty}}^2 \left(1 + \log \frac{1}{|\sin(\theta - \theta')|}\right). \tag{4-12}$$

This implies a similar estimate for  $\phi_2$ , because  $\phi_2 = \phi_2^{\#} + (I - K)^{-1}K\phi_2^{\#}$ . We summarize those estimates:

PROPOSITION 4.1. For  $\varepsilon > 0$  small enough, the fundamental solution  $\phi$  of (4–7) defined by (4–9) admits the representation  $\phi = \phi_0 + \phi_1 + \phi_2$  with

$$\phi_0 = E(x, \theta, -\infty)\delta(\theta - \theta') \int_0^{\tau_+(x', v(\theta'))} \delta(x - x' - tv(\theta')) dt,$$

$$\phi_1 = \chi E(y, \theta', -\infty) \frac{k(y, \theta', \theta)}{|\sin(\theta - \theta')|} E(y, \theta, \infty),$$

$$0 \le \phi_2 \le C ||k||_{L^\infty}^2 \left( 1 + \log \frac{1}{|\sin(\theta - \theta')|} \right),$$

where y,  $\chi$  are as in (4–11) and  $C = C(\operatorname{diam} X)$ .

Note that  $\phi_1$  is not a delta function but it is still singular with singularity at  $v(\theta) = v(\theta')$  (forward scattering) and  $v(\theta) = -v(\theta')$  (back-scattering). This singularity is integrable however, in fact it is easy to see that  $\int_{\Gamma_+} \phi_1 d\xi \le \int \sigma_p(x'+tv(\theta')) dt$ . The term  $\phi_2$  is also singular at  $v(\theta) = \pm v(\theta')$  with a weaker, logarithmic singularity. Therefore, we can still distinguish between the singularities of the three terms as in the case  $n \ge 3$ . This analysis however can give us information about k only near forward and backward directions. It is interesting to see whether by studying subsequent lower order terms we can recover all derivatives of  $k(x,\theta',\theta)$  at  $\theta=\theta'$  and  $\theta=\theta'+\pi$ . If so, this would allow us to recover collision kernels analytic in  $\theta$ ,  $\theta'$  (actually, analytic in  $\theta-\theta'$  would be enough) and to approximate k near  $\theta=\theta'$  and  $\theta=\theta'+\pi$  for smooth k. This would not require smallness assumptions on k but we still have to be sure that the direct problem is solvable.

We are ready now to sketch the proof of Theorem 4.1. Fix  $\Sigma > 0$  and assume that we have two pairs  $(\sigma_a, k)$ ,  $(\tilde{\sigma}_a, \tilde{k})$  in  $\mathcal{U}_{\Sigma,\varepsilon}$  with the same albedo operator and  $\sigma_a$ ,  $\tilde{\sigma}_a$  depending on x only. Denote by  $\phi_j$  and  $\tilde{\phi}_j$ , j = 0, 1, 2, the corresponding components of the fundamental solutions  $\phi$  and  $\tilde{\phi}$  as in Proposition 4.1. Then

$$\phi_0 + \phi_1 + \phi_2 = \tilde{\phi}_0 + \tilde{\phi}_1 + \tilde{\phi}_2 \quad \text{for } (x, \theta) \in \Gamma_+.$$
 (4-13)

As in Section 3, the most singular terms must agree; therefore,

$$\phi_0(x', \theta', x, \theta) = \tilde{\phi}_0(x', \theta', x, \theta) \qquad \text{for } (x', \theta') \in \Gamma_-, (x, \theta) \in \Gamma_+. \tag{4-14}$$

Thus the X-ray transform of  $\sigma_a(x)$  and  $\tilde{\sigma}_a(x)$  coincide. Therefore  $\sigma_a(x) = \tilde{\sigma}_a(x)$ . Next,

$$\chi E(y, \theta', -\infty) \frac{k(y, \theta', \theta) - \tilde{k}(y, \theta', \theta)}{|\sin(\theta - \theta')|} E(y, \theta, \infty) = \tilde{\phi}_2 - \phi_2 \quad \text{on } \Gamma_- \times \Gamma_+,$$

where y,  $\chi$  are as in (4–11). This, together with (4–13) and (4–14), leads to the inequality

$$\begin{split} \chi |k(y,\theta',\theta) - \tilde{k}(y,\theta',\theta)| &\leq C \big| \sin(\theta - \theta') \big| \big| \tilde{\phi}_1 - \phi_1 \big| \\ &= C \big| \sin(\theta - \theta') \big| \big| \tilde{\phi}_2 - \phi_2 \big| \qquad \text{on } \Gamma_- \times \Gamma_+, \quad (4\text{--}15) \end{split}$$

where  $C = e^{2d\Sigma}$ , d = diam D. The rest of the proof is based on the estimate

$$\underset{\Gamma_{-} \times \Gamma_{+}}{\operatorname{ess \, sup}} \left| \sin(\theta - \theta') \right| \left| \tilde{\phi}_{2} - \phi_{2} \right| \leq C \varepsilon \|k - \tilde{k}\|_{L^{\infty}(X \times S^{1} \times S^{1})}$$

$$(4-16)$$

with C > 0 depending on diam D only. An essential role in its proof is played by Lemma 4.1. Observe that in particular, the factor  $\sin(\theta - \theta')$  cancels the weaker logarithmic singularity of  $\phi_2$  and  $\tilde{\phi}_2$ . Then by (4–15), (4–16),

$$||k - \tilde{k}||_{L^{\infty}(X \times S^1 \times S^1)} \le C\varepsilon ||k - \tilde{k}||_{L^{\infty}(X \times S^1 \times S^1)},$$

and for  $\varepsilon > 0$  small enough this implies  $k = \tilde{k}$ .

Sketch of the proof of Theorem 4.2. According to Proposition 4.1, for the distribution kernel  $\alpha(x,\theta,x',\theta')$  of  $\mathcal{A}$  we have the representation (4–3), where a is as in (4–4), while  $b=(\phi_1+\phi_2)|_{(x,\theta)\in\Gamma_+}$ . Then b is a function and an elementary calculation show that  $b\in L^\infty(\Gamma_+,L^1(\Gamma_-))$ . Proposition 4.1 shows that  $b\sin(\theta-\theta')\in L^\infty$ . Assume that we have two pairs of continuous functions  $(\sigma_a,k)$  and  $(\tilde{\sigma}_a,\tilde{k})$  in  $\mathcal{V}^s_{\Sigma,\varepsilon}$  with albedo operators  $\mathcal{A}$  and  $\tilde{\mathcal{A}}$ , respectively. In what follows we will denote the quantities  $a,b,\alpha$ , etc., related to the second pair by putting a tilde sign over it. Also, we will use the notation  $\Delta\mathcal{A}$  to denote the difference  $\Delta\mathcal{A}=\mathcal{A}-\tilde{\mathcal{A}}$ , and similarly  $\Delta a=a-\tilde{a}$ , etc. In this paper,  $\Delta$  never stands for the Laplacian.

Let  $\delta_1$  and  $\delta_2$  be as in (4–5). According to Proposition 4.1,  $\delta_2 < \infty$ . Our goal is to estimate  $\Delta \sigma_a$  and  $\Delta k$  in terms of  $\delta_1$  and  $\delta_2$ . Observe that, as in the uniqueness proof,  $\Delta e^{-a}$  and  $\Delta(\alpha_1 + \alpha_2)$  can be recovered from  $\Delta \alpha$  by separating the most singular part of  $\Delta \alpha$  from the rest. Therefore,  $\delta_1$  measures the magnitude of the singular part of  $\Delta \alpha$ , while  $\delta_2$  measures the magnitude of the regular part.

We start by estimating  $\Delta \sigma_a$ . By a result of Mukhometov [Mu],

$$\|\Delta\sigma_a\|_{L^2} \le C\delta_1. \tag{4-17}$$

Next we estimate  $\Delta k$  in terms of  $\delta_1$  and  $\delta_2$ . Set

$$E_1(y, \theta, \theta') = E(y, \theta', -\infty)E(y, \theta, \infty).$$

By (4-11), 
$$\sin |\theta - \theta'| \alpha_1(x, \theta, x', \theta') = \chi(E_1 k)(y, \theta', \theta)$$
 with y,  $\chi$  as in (4-11).

Our starting point is the relation  $\Delta(E_1k) = k\Delta E_1 + \tilde{E}_1\Delta k$ . Note first that  $|\Delta E_1| \leq 2d|\Delta\sigma_a| \leq 2d\delta_1'$ , where  $\delta_1' = ||\Delta\sigma_a||_{L^{\infty}}$ . Hence, on supp  $\chi$ ,

$$\tilde{E}_{1}|\Delta k|(y,\theta,\theta') \leq |\Delta(E_{1}k)| + k|\Delta E_{1}| 
\leq |\Delta(\alpha_{1} + \alpha_{2})\sin|\theta - \theta'| - \Delta\alpha_{2}\sin|\theta - \theta'|| + C\delta'_{1} 
\leq \delta_{2} + |\Delta\alpha_{2}|\sin|\theta - \theta'| + C\delta'_{1}.$$
(4-18)

Therefore.

$$\|\Delta k\|_{L^{\infty}} \le C(\delta_1' + \delta_2 + \|\Delta \alpha_2 \sin |\theta - \theta'|\|_{L^{\infty}}) \tag{4-19}$$

The next step is to prove an estimate similar to (4–16) for the last term in the right-hand side above. Recall that in (4–16), we have  $\sigma_a = \tilde{\sigma}_a$ , while here we only know how to estimate the difference  $\Delta \sigma_a = \sigma_a - \tilde{\sigma}_a$ . Nevertheless, we can proceed along similar lines in order to get

$$|\Delta \alpha_2 \sin |\theta - \theta'||_{L^{\infty}} \le C\varepsilon (||\Delta k||_{L^{\infty}} + \delta_1'). \tag{4-20}$$

Therefore, for  $\varepsilon > 0$  small enough, (4–19) implies the stability estimate

$$\|\Delta k\|_{L^{\infty}} \le C\left(\delta_1' + \delta_2\right). \tag{4-21}$$

Estimate (4–21) is the base of our stability estimate. Using an interpolation inequality, we conclude that, for any fixed  $0 \le \mu \le 1 - 1/s$ ,

$$\|\Delta \sigma_a\|_{1+s\mu} \le C(\Sigma) \|\Delta \sigma_a\|_{L^2}^{1-1/s-\mu} \quad \text{for } \sigma_a, \tilde{\sigma}_a \in \mathcal{V}_{\Sigma,\varepsilon}^s.$$

By a standard Sobolev embedding theorem and (4–17),

$$\delta_1' = \|\Delta \sigma_a\|_{L^{\infty}} \le C(\Sigma) \|\Delta \sigma_a\|_{L^2}^{1-1/s-\mu} \le C'(\Sigma) \delta_1^{1-1/s-\mu}. \tag{4-22}$$

Therefore, (4–21) yields

$$\|\Delta k\|_{L^{\infty}} \le C\left(\delta_1^{1-1/s-\mu} + \delta_2\right). \tag{4-23}$$

Estimates (4-22) and (4-23) complete the proof of Theorem 4.2.

#### 5. Open Problems

Our choice of open problems is subjective and mainly reflects personal taste.

Uniqueness for  $\sigma_a$  depending on both x and v. As we have demonstrated in the previous sections, if k=0, then  $\sigma_a(x,v)$  cannot be recovered from  $\mathcal{A}$  (or from the scattering operator) because the line integrals  $\int \sigma_a(x+tv,v) dt$  do not recover  $\sigma_a$ . Suppose however that  $\sigma_p = \sigma_a$ . Then the absorption is due only to the fact that particles may change velocity and each such event is interpreted as a particle instantly moving from the point (x,v') of the phase space (therefore absorption at (x,v')) to the point (x,v). Then the counter example above does not work. Can we recover  $\sigma_a(x,v)$  in this case? If yes, then recovery of k goes along the same lines as above. More generally, one can assume

that  $\sigma_a(x, v) = \sigma_p(x, v) + a(x)$ . Even more generally, the counter example above cannot be generalized in an obvious way if k > 0 in X. Is this condition alone enough for uniqueness?

Relaxing the smallness condition in the two-dimensional case. It would be interesting to prove uniqueness in the 2D case without smallness assumption on k. Some conditions on  $\sigma_a$  and k are needed even for the direct problem; see, for example, (3–1) and (3–3) — the first one does require k to be small (but with an explicit bound in general much larger than the one needed for the inverse problem) while the second one does not. As mentioned in Section 4, one can try to recover k(x, v', v) near  $v = \pm v'$  at infinite order by studying the singularities of the kernel of  $\mathcal{A}$  which solves the problem for k analytic in v and v'. It would also be interesting to see whether one could recover the singularities of k from boundary measurements, at least if we assume that they are of jump type across some curve.

Stability for  $n \geq 3$ . Stability estimates in dimensions  $n \geq 3$  have been proven by Romanov (see [R1], [R2], [R3], for example) and by Wang [W], under additional assumptions that k depend on fewer variables. In the general situation studied in Section 3, we know of no stability estimates, even for small k as in Theorem 4.2 (where n=2). We believe that such an estimate should be possible to derive following the proof of Theorem 3.1. This is done in [W] under the additional assumption that k = k(v', v).

Alternative recovery method for large  $\sigma_a$  and k. We do not impose smallness assumptions on the coefficients in dimensions  $n \geq 3$ , and our method gives in fact an explicit solution of the inverse problem which in particular implies a reconstruction method based on taking certain limits near the singularity of  $\alpha$ . However, for large  $\sigma_a$ , the amplitude of the most singular part  $\alpha_0$  is exponentially small for large  $\sigma_a$ . For all practical purposes, measuring the leading singularity is hard or impossible in this case. Therefore, it would be important to develop a method for relatively large  $\sigma_a$  and k that does not rely on measuring the singularities of  $\alpha$ . One possible way is to study the diffusion limit (replacing  $\sigma_a$  and k by  $\lambda \sigma_a$  and  $\lambda k$  and taking  $\lambda \to \infty$ ) and an associated inverse problem.

#### References

- [AB] Yu. E. Anikonov, B. A. Bubnov, Inverse problems of transport theory, Soviet Math. Doklady 37 (1988), 497–499.
- [AE] P. Arianfar and H. Emamirad, Relation between scattering and albedo operators in linear transport theory, Transport Theory and Stat. Physics, 23(4) (1994), 517–531.
- [B] A. Bondarenko, The structure of the fundamental solution of the time-independent transport equation, J. Math. Anal. Appl. 221:2 (1998), 430–451.
- [Ce1] M. Cessenat, Théorèmes de trace L<sup>p</sup> pour des espaces de fonctions de la neutronique, C. R. Acad. Sci. Paris, Série I 299 (1984), 831–834.

- [Ce2] M. Cessenat, Théorèmes de trace pour des espaces de fonctions de la neutronique,
   C. R. Acad. Sci. Paris, Série I 300 (1985), 89-92.
- [CMS] M. Chabi, M. Mokhtar-Kharroubi and P. Stefanov, Scattering theory with two L<sup>1</sup> spaces: application to transport equations with obstacles, Ann. Fac. Sci. Toulouse 6(3) (1997), 511–523.
- [CSt1] M. Choulli and P. Stefanov, Scattering inverse pour l'équation du transport et relations entre les opérateurs de scattering et d'albédo, C. R. Acad. Sci. Paris 320 (1995), 947–952.
- [CSt2] M. Choulli and P. Stefanov, Inverse scattering and inverse boundary value problems for the linear Boltzmann equation, Comm. P.D.E. 21 (1996), 763–785.
- [CSt3] M. Choulli and P. Stefanov, Reconstruction of the coefficients of the stationary transport equation from boundary measurements, Inverse Prob. 12 (1996), L19–L23.
- [CSt4] M. Choulli and P. Stefanov, An inverse boundary value problem for the stationary transport equation, Osaka J. Math. 36(1) (1999), 87–104.
- [CZ] M. Choulli, A. Zeghal, Laplace transform approach for an inverse problem, Transport Theory and Stat. Phys., 24(9) (1995), 1353–1367.
- [DL] R. Dautray and J.-L. Lions, Analyse Mathématique et Calcul Numérique pour les Sciences et les Techniques, vol. 9, Masson, Paris, 1988.
- [E] H. Emamirad, On the Lax and Phillips scattering theory for transport equation,
   J. Funct. Anal. 62 (1985), 276–303.
- [EP] H. Emamirad and V.Protopopescu, Relationship between the albedo and scattering operators for the Boltzmann equation with semi-transparent boundary conditions, Math. Meth. Appl. Sci., to appear.
- [Hej] J. Hejtmanek, Scattering theory of the linear Boltzmann operator, Commun. Math. Physics 43 (1975), 109–120.
- [H] S. Helgason, The Radon Transform, Birkhäuser, Boston, Basel, 1980.
- [L1] E. W. Larsen, Solution of multidimensional inverse transport problems, J. Math. Phys. 25(1) (1984), 131–135.
- [L2] E. Larsen, Solution of three-dimensional inverse transport problems, Transport Theory Statist. Phys. 17:2–3 (1988), 147–167.
- [MC] N. J. McCormick, Recent developments in inverse scattering transport methods, Trans. Theory and Stat. Phys. 13 (1984), 15–28.
- [Mo] M. Mokhtar-Kharroubi, Limiting absorption principle and wave operators on  $L^1(\mu)$  spaces, Applications to Transport Theory, J. Funct. Anal. **115** (1993), 119-145.
- [Mu] R. Mukhometov, On the problem of integral geometry, Math. Problems of Geophysics, Akad. Nauk SSSR, Sibirs. Otdel., Vychisl. Tsentr, Novosibirsk, 6:2 (1975), 212–242 (in Russian).
- [PV] A. I. Prilepko, N. P. Volkov, Inverse problems for determining the parameters of nonstationary kinetic transport equation from additional information on the traces of the unknown function, Differentsialnye Uravneniya 24 (1988), 136–146.
- [P1] V. Protopopescu, On the scattering matrix for the linear Boltzmann equation, Rev. Roum. Phys. 21 (1976), 991–994.

- [P2] V. Protopopescu, Relation entre les opérateurs d'albédo et de Scattering avec des conditions aux frontières non-transparentes, C. R. Acad. Sci. Paris, Série I 318 (1994), 83–86.
- [RS] M. Reed and B. Simon, Methods of Modern Mathematical Physics, Vol. 3, Academic Press, New York, 1979.
- [R1] V.G. Romanov, Estimation of stability in the problem of determining the attenuation coefficient and the scattering indicatrix for the transport equation, Sibirsk. Mat. Zh. 37:2 (1996), 361–377, iii; translation in Siberian Math. J. 37:2 (1996), 308–324.
- [R2] V.G. Romanov, Stability estimates in the three-dimensional inverse problem for the transport equation, J. Inverse Ill-Posed Probl. 5:5 (1997), 463–475
- [R3] V. Romanov, A stability theorem in the problem of the joint determination of the attenuation coefficient and the scattering indicatrix for the stationary transport equation, Mat. Tr. 1:1 (1998), 78–115.
- [Si] B. Simon, Existence of the scattering matrix for linearized Boltzmann equation, Commun. Math. Physics 41 (1975), 99–108.
- [St] P. Stefanov, Spectral and scattering theory for the linear Boltzmann equation, Math. Nachr. 137 (1988), 63–77.
- [StU] P. Stefanov and G. Uhlmann, Optical Tomography in two dimensions, to appear in Methods Appl. Anal.
- [SyU] J. Sylvester and G. Uhlmann, The Dirichlet to Neumann map and applications, in: Inverse Problems in Partial Differential Equations, SIAM Proceedings Series List (1990), 101–139, ed. by D. Colton, R. Ewing and R. Rundell.
- [T] A. Tamasan, An inverse boundary value problem in two dimensional transport, Inverse Problems 18 (2002), 209-219.
- [U] T. Umeda, Scattering and spectral theory for the linear Boltzmann operator,
   J. Math. Kyoto Univ. 24 (1984), 208–218.
- [Vi] I. Vidav, Existence and uniqueness of non-negative eigenfunctions of the Boltzmann operator, J. Math. Anal. Appl. 22 (1968), 144-155.
- [V1] J. Voigt, On the existence of the scattering operator for the linear Boltzmann equation, J. Math. Anal. Appl. **58** (1977), 541–558.
- [V2] J. Voigt, Spectral properties of the neutron transport equation, J. Math. Anal. Appl. 106 (1985), 140–153.
- [W] J. Wang, Stability estimates of an inverse problem for the stationary transport equation, Ann. Inst. H. Poincaré Phys. Théor. 70:5 (1999), 473–495.

PLAMEN STEFANOV
DEPARTMENT OF MATHEMATICS
PURDUE UNIVERSITY
WEST LAFAYETTE, IN 47907
UNITED STATES
stefanov@math.purdue.edu