An Inverse Scattering Problem for Dirac Equations with Time-Dependent Electromagnetic Potentials

Dedicated to the memory of Professor Nobuhisa Iwasaki

By

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Abstract

We consider an inverse scattering problem, by using a time-dependent method, for the Dirac equation with a time-dependent electromagnetic field. The Fourier transform of the field is reconstructed from the scattering operator on a Lorents invariant set

 $(0.1) D:=\{(\tau,\xi)\in \mathbb{R}\times\mathbb{R}^3; |\tau|< c|\xi|\}$

in the dual space of the space-time. As corollaries of this result, we can reconstruct the electromagnetic field completely if it is a finite sum of fields each of which is a time-independent one by a suitable Lorentz transform, and we can also determine the field uniquely if the fields satisfies some exponential decay condition. Our assumptions and results are independent of a choice of inertial frames.

§1. Introduction

In this paper we consider an inverse scattering problem for the Dirac equation with a time-dependent electromagnetic field.

If the field is time-independent and short range, then it can be completely reconstructed from the scattering operator (see [Is], [It], [J]). This means that the field can be determined by scattering experiments in the inertial frame in which the field is time-independent. Because of the relativistic invariance of the Dirac equation, it follows that the field can be also determined by scattering experiments in any inertial frame, the field is in general time-dependent. Thus, it is important to study the Dirac equation with time-dependent fields for the

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inverse scattering problem.

We proceed our argument by fixing an inertial frame; however, our assumptions and results are stated in a way independent of a choice of inertial frames (see Section 5). The variable t denotes the time and x is the space-variable.

We begin with some explanation for our notations. We denote by $\langle a,b \rangle_{R^d}$ the usual inner product of a and b in \mathbb{R}^d and may write $a \cdot b$ or $\langle a,b \rangle$ for simplicity. Moreover, the usual norm of \mathbb{R}^d is denoted by the same symbol $|\cdot|$ for any d if no confusion occures. We also use the symbol $\langle T,u \rangle$ in place of T(u) for a distribution T and a test function u.

The Dirac equation with an electromagnetic potential

$$A = (A^0, A^1, A^2, A^3) : \mathbb{R}_t \times \mathbb{R}_x^3 \longrightarrow \mathbb{R}^4$$

is given by

(1.1)
$$i\frac{d}{dt}\Psi(t) = H_A(t)\Psi(t), \quad \Psi(t) \in \mathcal{H} := L^2(\mathbb{R}^3_x; \mathbb{C}^4),$$

$$H_A(t) = c \sum_{j=1}^{3} \alpha_j (D_j - A^j(t,x)) + \alpha_4 m c^2 - A^0(t,x) I_4,$$

where c > 0 is the velocity of light, $m \ge 0$ the rest mass of the particle, $D_j = -i\partial/\partial x_j$, and α_j 's are 4×4 Hermitian matrices with the following properties:

$$\alpha_j \alpha_k + \alpha_k \alpha_j = 2\delta_{jk} I_4, \quad 1 \leq j,k \leq 4,$$

where δ_{jk} is the Kronecker symbol and I_n is the $n \times n$ identity matrix.

Let $L \neq \{0\}$ be a subspace of \mathbb{R}^4 . Then we denote by X_L the orthogonal projection of $X = (t,x) \in \mathbb{R}^4$ onto L and define a class of potentials:

$$S(L):=\{A\in \mathscr{B}^1(\mathbb{R}^4,\mathbb{R}^4);\int_0^\infty g_A^L(r)\,dr<\infty\}$$

where $g_A^L(r) := \sup_{|X_l| \ge r} |A(X)|$ and $\mathscr{B}^1(\mathbb{R}^4, \mathbb{R}^4)$ is the space of bounded $C^1(\mathbb{R}^4, \mathbb{R}^4)$ -functions with bounded derivatives. If $A \in \mathscr{B}^1(\mathbb{R}^4, \mathbb{R}^4)$ satisfies the short range condition with respect to X_L -variable:

$$|A(X)| \leq K(1+|X_L|)^{-\rho}$$
 on \mathbb{R}^4

for some K > 0 and some $\rho > 1$, then A belongs to S(L).

We also say that A belongs to S if and only if A is decomposed as

$$A = \sum_{j=1}^{N} A_j, \quad A_j \in S(L_j),$$

for some N and for some subspaces L_j , $1 \le j \le N$.

If A belongs to S, the Dirac equation (1.1) has a unique unitary propagator $U_A(t,s)$, $s,t \in \mathbb{R}$:

$$i\frac{d}{dt}U_A(t,s) = H_A(t)U_A(t,s), \quad U_A(s,s) = I_A(t)U_A(t,s),$$

Moreover, we have the following.

Proposition 1.1. Let $A \in S$. Then the wave operators

$$W_A^{\pm}(s) := s - \lim_{t \to \pm \infty} U_A(s,t) e^{-t(t-s)H_0}$$

exist for each $s \in \mathbb{R}$, where H_0 is the free Dirac operator.

$$H_0 = c \sum_{j=1}^3 \alpha_j D_j + \alpha_4 m c^2.$$

Remark. The free Dirac operator H_0 is a self-adjoint operator with domain $D(H_0) = H^1(\mathbb{R}^3, \mathbb{C}^4)$, the Sobolev space of order 1, and $U_0(t,s) = e^{-i(t-s)H_0}$.

The scattering operator is defined by

$$S_A(s) := W_A^+(s)^* W_A^-(s),$$

for each $s \in \mathbb{R}$. If some stronger condition is imposed on A, the scattering operator is unitary in \mathcal{H} . But, it is not necessarily unitary under the assumption $A \in S$.

The following useful relation follows immediately from definition:

(1.2)
$$S_A(s) = e^{-isH_0}S_A(0)e^{isH_0}, s \in \mathbb{R}$$

Thanks to this relation, we know $S_A(s)$ for all $s \in \mathbb{R}$ if $S_A(s_0)$ is given for some s_0 . The electromagnetic field F_A , determined by A, is defined by

$$F_A = (F_A^{jk})_{0 \le j < k \le 3} = \left(\frac{\partial A^k}{\partial x_j} - \frac{\partial A^j}{\partial x_k}\right)_{0 \le j < k \le 3} : \mathbb{R}^4 \longrightarrow \mathbb{R}^6,$$

where $x_0 = t$.

It should be recalled that the potential is not uniquely determined by the field and that it is not the potential but the field that can be an observable quantity. The following theorem shows that the scattering operator is determined by the field not by the potential.

Theorem 1.2. Let $A_{(1)}$ and $A_{(2)}$ belong to S and suppose that

$$A_{(2)} - A_{(1)} = \sum_{j=1}^{N} A_j, \quad A_j \in S(L_j), \quad dimL_j \ge 2, \ 1 \le j \le N,$$

and that $F_{A_{(1)}} = F_{A_{(2)}}$. Then $S_{A_{(1)}}(s) = S_{A_{(2)}}(s)$ for all $s \in \mathbb{R}$.

We next consider the inverse problem.

For a subspace $L \neq \{0\}$ of \mathbb{R}^4 the class $\widetilde{S}(L)$ of electromagnetic fields is defined by

$$\widetilde{S}(L) := \{F \in \mathscr{B}^0(\mathbb{R}^4, \mathbb{R}^6); \int_0^\infty g_F^L(r) dr < \infty\},\$$

where $g_F^L(r) := \sup_{\substack{|X_{cl}| \geq r \\ \text{functions from } \mathbb{R}^4 \text{ to } \mathbb{R}^6} |F(X)|$ and $\mathscr{B}^0(\mathbb{R}^4, \mathbb{R}^6)$ is the space of bounded continuous

We denote by $\mathcal{Z}=(\tau,\xi)\in\mathbb{R}\times\mathbb{R}^3$ the dual variables of X=(t,x) and define an open set D in \mathbb{R}^4 by

(1.3)
$$D := \{ (\tau, \xi) \in \mathbb{R}^4; |\tau| < c |\xi| \}$$

We denote by $\mathscr{S}'(\mathbb{R}^4; \mathbb{C}^6)$ the space of \mathbb{C}^6 -valued tempered distributions and by $\widehat{F} \in \mathscr{S}'(\mathbb{R}^4; \mathbb{C}^6)$ the Fourier transform of $F \in \mathscr{S}'(\mathbb{R}^4; \mathbb{C}^6)$:

$$\widehat{F}(\tau,\xi) = (2\pi)^{-2} \iint e^{-it\tau - ix \cdot \xi} F(t,x) dt dx.$$

We denote by $\widehat{F}|_{\mathscr{Q}}$ the restriction of \widehat{F} to an open set $\mathcal{Q}, \widehat{F}|_{\mathscr{Q}} \in \mathscr{D}'(\mathcal{Q})$.

Theorem 1.3. Suppose that $A = \sum_{j=1}^{N} A_j$, with $A_j \in S(L_j)$ and that $F_A = \sum_{k=1}^{N'} F_k$ with $F_k \in \widetilde{S}(L'_k)$. Then $\widehat{F_A}|_{D\setminus\Sigma}$ can be reconstructed from $S_A(0)$, where

(1.4)
$$\sum := \left(\bigcup_{\substack{1 \le j \le N \\ \dim L_j = 1}} L_j\right) \cup \left(\bigcup_{\substack{1 \le k \le N \\ \dim L_k = 1}} L_k'\right).$$

Remark. The decompositions of A and F_A are not unique and the exceptional set Σ depends on the decompositions. Let \mathscr{A} be the set of all pairs of decompositions of A and F_A , and denote by Σ_a the Σ in (1.4) for $a \in \mathscr{A}$. Then Σ in Theorem 1.3 may be replaced by the smaller set

$$\widetilde{\Sigma} := \bigcap_{a \in \mathscr{A}} \Sigma_a$$

Theorem 1.3 tells us nothing about \widehat{F}_A on D^c , the complement of D. If \widehat{F}_A

on D^c does not contributes to the scattering operator, this result would be satisfactory. However, the following Proposition shows that this is not the case.

Proposition 1.4 Let c = m = 1 and let

$$D_1 := \{(\tau, \xi) \in \mathbb{R} \times \mathbb{R}^3; \tau > \sqrt{|\xi|^2 + 4}\}.$$

Then for any $\phi \in \mathscr{S}(\mathbb{R}_t \times \mathbb{R}^3_x; \mathbb{R})$ such that $supp \widehat{\phi} \cap D_1 \neq \emptyset$, the scattering operator $S_{\lambda A}(0)$ for (1.1) with the potential $\lambda A = (\lambda \phi, 0, 0, 0)$ satisfies $S_{\lambda A}(0) \neq I$ for any small $\lambda \neq 0$.

Remark. The author does not know whether the field F_A can be reconstructed completely from the scattering operator if the support of \widehat{F}_A has an intersection with D^c , while on the other hand an affirmative answer is known in the case of Schrödinger equations [We1]. The velocity of a relativistic particle cannot exceed that of light, whereas a nonrelativistic particle can have any speed. This is one of the greatest differences between a particle obeying a Dirac equation and one obeying a Schrödinger equation. The following intuitive argument shows that the field F with supp $\widehat{F} \subset D' := \{(\tau, \xi); |\tau| > c|\xi|\}$ may be peculiar one from the point of view of the relativity. We write formaly

$$F(t,x) = (2\pi)^{-2} \int_{D'} F_{\tau,\xi}(t,x) \, d\tau d\xi,$$

where $F_{\tau,\xi}(t,x) := e^{it\tau + ix \cdot \xi} \widehat{F}(\tau,\xi)$ satisfies the wave equation

$$\frac{1}{\tau^2 |\xi|^{-2}} \frac{\partial^2}{\partial t^2} F_{\tau,\xi}(t,x) = \Delta_x F_{\tau,\xi}(t,x).$$

This implies that each component $F_{\tau,\xi}$ of the field F propagates with velocity $|\tau||\xi|^{-1} > c$, which contradicts the relativity.

We can determine the field completely, if we impose some conditions on the field, as corollaries of Theorem 1.3.

Theorem 1.5. Suppose (i) and (ii) in Theorem 1.3. Moreover, we assume $\sum \cap D = \emptyset$ and

(1.5)
$$(supp \widehat{F}_A \setminus \{0\}) \cap D^c = \emptyset.$$

Then F_A can be completely reconstructed from $S_A(0)$.

Using this theorem we can treat time-independent potentials.

Corollary 1.6. Suppose (i) and (ii) in Theorem 1.3 with $\Sigma \cap D = \emptyset$. Moreover, suppose that for each $k=1,\dots,N'$ there exists $V_k \in T := \{X = (t,x) \in \mathbb{R}^4; c|t| > |x|\}$ such that

(1.6)
$$F_k(sV_k+X) = F_k(X), \quad s \in \mathbb{R}, \quad X \in \mathbb{R}^4.$$

Then F_A can be reconstructed completely from $S_A(0)$.

Remark 1. In Section 5 we will show that each F_k satisfying (1.6) may be made time-independent if a suitable inertial frame is chosen.

Remark 2. If $A \in S$ is independent of t and satisfies

(1.7)
$$|A(x)| + |F_A(x)| \le K(1+|x|)^{-\rho} \text{ on } \mathbb{R}^3$$

for some constants K>0 and $\rho>1$. Then Corollary 1.6 shows that $F_A(x)$ can be completely reconstructed from the scattering operator $S_A(0)$. This has been known under different conditions in [It], [J] and [Is]. Roughly speaking, the decay rate of the potential is supposed to be $\rho>3$ in [It], $\rho>3/2$ in [J] and $\rho>$ 2 in [Is]. However, the decay condition on the field is not imposed in [J], and the magnetic field is not treated in [Is].

The next theorem shows that the field is uniquely determined by the scattering operator under some exponential decay condition.

Theorem 1.7. Let $A \in S$. (1) Suppose there exists $V \in S^3$, S^3 being the unit sphere in \mathbb{R}^4 , such that

(1.8) $|F_A(X)| \le K e^{-\delta|\langle V, X \rangle_{R^4}} \text{ on } \mathbb{R}^4,$

for some constants K>0 and $\delta>0$.

Then $\widehat{F_A}|_{R^{N_L}}$ is determined by $S_A(0)$, where L is the one-dimensional subspace spanned by V.

(2) Suppose, in addition to the assumption of (1), that there exist $V \in S^3$ not parallel to V and a bounded function g satisfying $g(t)(1+t)^{-1} \in L^1((0,\infty))$ such that

(1.9)
$$|F_A(X)| \le g(|\langle V', X \rangle_{R^4}|) e^{-\delta|\langle V, X \rangle_{R^4}|} \text{ on } \mathbb{R}^4.$$

Then F_A is uniquely determined by $S_A(0)$.

Remark. Roughly speaking, (2) implies that if F_A satisfies

$$|F_A(X)| \leq K(1+|\langle V_1,X\rangle_{R'}|)^{-\rho} e^{-\delta|\langle V_2,X\rangle_{R'}|}$$
 on \mathbb{R}^4

for some linearly independent $V_1, V_2 \in \mathbb{R}^4 \cap S^3$ and for some K > 0, $\rho > 0$ and $\delta > 0$, then F_A is completely determined by the scattering operator.

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It is an important problem in physics as well as in mathematics to know whether the external field can be determined from the scattering operator. In the case of Schrödinger operators with time-independent potentials, it is known, since Faddeev [F], that the potential can be reconstructed from the high-energy behavior of the scattering matrices (see, for example, [Is-K], [Ne], [Sa], [Wa], [Ni]). The proofs are based on a stationary representation of the scattering matrices and some resolvent estimates at the high-energy range. Using a similar stationary method, the author [It] has proved that the electromagnetic filed can be reconstructed from the high-energy behavior of the scattering matrices of the Dirac operator with a time-independent potential. In [Is] Isozaki has obtained a similar result as well as an uniqueness result for the fixed energy problem by a different method.

On the other hand, Enss and Weder [E-We] have found a new method, a time-dependent method (a geometric method), to reconstruct the potential from the high-energy asymptotics of the scattering operator in the case of Schrödinger operators without magnetic fields. Since their method is simple, it can be applicable to many cases. Recently, Arians [A1] has applied their method to reconstruct the electromagnetic field for the Schrödinger operator with a time-independent electromagnetic potential. See also [A2, We2, We3]. On the other hand, Weder [We1] has shown that the potential can be completely reconstructed from the scattering operator for the Schrödinger equations with a time-dependent potential. For Dirac operators with a time-independent potential Jung [J] has reconstructed the electromagnetic field by using the geometric method. Some of his proofs are applicable to the case of time-dependent potentials (Proposition 2.2).

At the end of this section we describe the plan of our paper. Section 2contains the proofs of Proposition 1.1, Theorems 1.2 and 1.3. Proposition 1.1 is proved by Cook's method with the help of estimates of the propagator for the free Dirac operator in the classically forbidden region. The proof of Theorem 1.2 is carried out by constructing a function u_1 satisfying $\nabla_x u_1 = A$ and decaying in some directions. The proof of Theorem 1.3 is based on Proposition 2.2, which gives high-energy asymptotics of the scattering operator. Proposition 2.2 is proved by using a time-dependent method, due to Enss and Weder, as in [J]. By means of Proposition 2.2 and Stoke's theorem we can prove that some line integrals (2.31) for the electromagnetic field are constructed from the scattering operator. If the magnetic field does not exist, the line integrals are the same ones as appeared in [St] and [R-S], in which the inverse scattering for wave equations with time-dependent potentials is studied. Roughly speaking, Theorem 1.3 is proved by constructing the field from the line integrals with the help of the Fourier transform. The proofs of Theorems 1.5 and 1.7, based on Theorem 1.3, are given in Section 3. Proposition 1.4 is proved by investigating the second term of the Dyson expansion of the scattering operator in Section 4.

Dirac equations are relativistic invariant. In Section 5 we will show that our assumptions and results are independent of the choice of the inertial frame.

§ 2. Proofs of Proposition 1.1, Theorems 1.2 and 1.3

In this section we will give proofs of Proposition 1.1, Theorem 1.2 and 1.3.

Proof of Proposition 1.1. Let
$$A = \sum_{j=1}^{N} A_j$$
, $A_j \in S(L_j)$ and define
 $K_A := \bigcup_{j=1}^{N} \{\xi \in \mathbb{R}^3; V^+(\xi) \in L_j^{\perp} \text{ or } V^-(\xi) \in L_j^{\perp}\}.$

where

$$V^{\pm}(\xi) := (1, \pm v(\xi)) = (1, \pm c^2 \xi / \sqrt{c^2 |\xi|^2 + m^2 c^4})$$

and L^{\perp} is the orthogonal complement of $L \subset \mathbb{R}^4$. Since K_A is a closed null set in \mathbb{R}^3 ,

$$G := \{ f \in \mathscr{S}(\mathbb{R}^3, \mathbb{C}^4); \, \widehat{f} \in C_0^{\infty}(\mathbb{R}^3 \setminus (K_A \cup \{0\}); \mathbb{C}^4) \}$$

is dense in \mathcal{H} , where \hat{f} is the Fourier transform of f. By Cook's method, Proposition 1.1 follows from it:

(2.1)
$$\int_{-\infty}^{\infty} ||A_j(t, \cdot)e^{-\iota t H_0} P_{\pm} f|| dt < +\infty, \quad j=1, \cdots N,$$

for all $f \in G$, where P_{\pm} is the orthogonal projection in \mathscr{H} onto the positive/negative energy subspace of H_0 ,

$$P_{\pm} := \frac{1}{2} \left(I \pm \frac{H_0}{|H_0|} \right),$$

and $\|\cdot\|$ denotes the norm of \mathcal{H} .

We have for $\delta > 0$

(2.2)
$$\|A_{j}(t, \cdot)e^{-itH_{0}}P_{\pm}f\|$$

$$\leq \|A_{j}(t, \cdot)\chi(|(t, x)_{L_{j}}| \geq \delta|t|)\|_{L^{\infty}}\|f\|$$

$$+ \|A_{j}(t, \cdot)\|_{L^{\infty}}\|\chi(|(t, x)_{L_{j}}| \leq \delta|t|)e^{-itH_{0}}P_{\pm}f\|$$

where $\chi(M(x))$ denotes the multiplication operator by the characteristic function of x satisfying the condition M(x) and $(t,x)_L$ is the orthogonal projection of $(t,x) \in \mathbb{R}^4$ onto L. Since $A_j \in S(L_j)$, the first term on the right is bounded by $g_{k'}(\delta|t|) ||f||$ and so is integrable with respect to $t \in \mathbb{R}$. Therefore it suffices to show that there exists a small $\delta > 0$ such that

(2.3)
$$\|\chi(|(t,x)_{L_{1}}| \leq \delta|t|) e^{-itH_{0}}P_{\pm}f\| \leq K_{n}(1+|t|)^{-n}$$

for any $n \ge 0$, where K_n is a constant independent of t. The symbol of H_0 .

$$h_0(\xi) := c \sum_{j=1}^{3} \alpha_j \xi_j + \alpha_4 m c^2,$$

has eigenvalues $\pm \sqrt{c^2 |\xi|^2 + m^2 c^4}$ and the associated eigenprojection is given by

$$P^{\pm}(\xi) = \frac{1}{2} \Big(I_4 \pm \frac{h_0(\xi)}{\sqrt{c^2 |\xi|^2 + m^2 c^4}} \Big).$$

Hence we have

$$(e^{-itH_0}P_{\pm}f)(x) = (2\pi)^{-3/2} \int e^{ig^{\pm}(t,x,\xi)} \widehat{f}_{\pm}(\xi) d\xi,$$

where $g^{\pm}(t,x,\xi) = \langle x,\xi \rangle \mp t \sqrt{c^2 |\xi|^2 + m^2 c^4}$ and $\hat{f}_{\pm}(\xi) = P^{\pm}(\xi) \hat{f}(\xi)$. We claim that there exists a constant K > 0 such that

$$(2.4) \qquad \qquad |\nabla_{\xi}g^{\pm}(t,x,\xi)| \ge K(|t|+|x|)$$

for all $(t, x, \xi) \in \Omega_j \times \text{supp } \widehat{f}$ if $\delta > 0$ is sufficiently small, where

$$\Omega_{j}:=\{(t,x)\in \mathbb{R}\times\mathbb{R}^{3}; |t|>1, |(t,x)_{L_{j}}|\leq\delta|t|\}.$$

Since $|v(\xi)| \leq c$ and

(2.5)
$$\nabla_{\xi}g^{\pm}(t,x,\xi) = x \mp tv(\xi),$$

it follows that

$$(2.6) |\nabla_{\xi}g^{\pm}| \ge |x| - c|t|.$$

On the other hand, for $(t, x, \xi) \in \Omega_j \times \text{supp } \widehat{f}$, we have

$$|x \mp tv(\xi)| = |(t,x) - tV^{\pm}(\xi)| \ge |(t,x)_{L_{j}} - tV^{\pm}(\xi)_{L_{j}}|$$
$$\ge |t|(|V^{\pm}(\xi)_{L_{j}}| - \delta) \ge (\delta_{1} - \delta)|t|,$$

where $\delta_1 := \inf_{\xi \in \operatorname{supp} \widehat{f}} |V^{\pm}(\widehat{\xi})_{L_j}| > 0$, since dist $(K_A, \operatorname{supp} \widehat{f}) > 0$. Thus, if $0 < \delta < \delta_1$, taking account of (2.6), we obtain (2.4). Therefore, it is easy to see that, on $\mathcal{Q}_j \times \operatorname{supp} \widehat{f}$, $\nabla_{\xi} g^{\pm} / |\nabla_{\xi} g^{\pm}|^2 \in C^{\infty}$ and

(2.7)
$$\left|\partial_{\xi}^{r}\left(\frac{\nabla_{\xi}g^{\pm}}{|\nabla_{\xi}g^{\pm}|^{2}}\right)\right| \leq K_{r}(|t|+|x|)^{-1}$$

for any γ . Thus, using integration by parts together with

$$ie^{ig^{\pm}} = rac{
abla^{\pm}g}{|
abla_{\xi}g^{\pm}|^2} \cdot
abla_{\xi}e^{ig^{\pm}},$$

we conclude that for any n there exists K_n such that

(2.8)
$$|(e^{-\iota t H_0} P_{\pm} f)(x)| \leq K_n (1+|t|+|x|)^{-n},$$

for all $(t,x) \in \Omega_j$. This implies (2.3) and completes the proof of Proposition 1.1. \Box

Remark. If L is a nontrivial subspace of \mathbb{R}^4 and if $J(\sigma)$ is a monotone decreasing function on $(0,\infty)$ with $\lim J(\sigma)=0$, then

(2.9)
$$s - \lim_{t \to \pm \infty} J(|(t, \cdot)_L|) e^{-tH_0} = 0.$$

This can be proved in the same way as above by replacing $A_j(t,x)$ by $J(|(t,x)_L|)$. This fact will be used in the proof of Lemma 2.1 below.

Proof of Theorem 1.2. Let
$$A := A_{(2)} - A_{(1)}$$
. Since $F_A = 0$,
$$u(X) := \int_0^1 \langle X, A(sX) \rangle_{R'} ds$$

is in $C^1(\mathbb{R}^4)$ and satisfies

 $(2.10) \nabla_X u = A \text{ on } \mathbb{R}^4.$

Let $V \in S^3 \setminus (\bigcup_{j=1}^{N} L_j^{\perp})$, where S^3 is the unit sphere in \mathbb{R}^4 . Then $V_{L_j} \neq 0$ for j=1, \dots, N . Using (2.10), we have for $0 < r_1 < r_2$

$$(2.11) |u(r_2V) - u(r_1V)| = \left| \int_{r_1}^{r_2} \frac{d}{dr} u(rV) dr \right| \\ \leq \int_{r_1}^{r_2} |\langle V, A(rV) \rangle_{R^2} |dr| \\ \leq \sum_{j=1}^N \int_{r_1}^{r_2} g_{A_j}^{L_j}(r|V_{L_j}|) dr.$$

(2.11) implies that the limit

$$u_{\infty}(V) := \lim_{r \to +\infty} u(rV)$$

exists. We next prove that the limit $u_{\infty}(V)$ is independent of $V \in S^3 \setminus (\bigcup_{j=1}^N L_j^{\perp})$. Let $V, V' \in S^3 \setminus (\bigcup_{j=1}^N L_j^{\perp})$. Since dim $L_j^{\perp} \leq 2, S^3 \setminus (\bigcup_{j=1}^N L_j^{\perp})$ is connected. Hence, we may prove $u_{\infty}(V) = u_{\infty}(V')$ for V, V' such that the segment between V and V' has no intersection with $\bigcup_{j=1}^N L_j^{\perp}$. Then there exists $\varepsilon_0 > 0$ such that

$$\inf_{\sigma\in[0,1]} \left| \sigma(V_{L_j} - V_{L_j}) + V_{L_j} \right| > \varepsilon_0, \quad j=1,\cdots,N.$$

Using (2.10) again, we estimate

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$$|u(rV) - u(rV')| \leq \int_0^1 |r(V' - V)| |A(r(\sigma(V' - V) + V))| d\sigma$$
$$\leq 2 \sum_{j=1}^N rg_{A_j}^{L_j}(\varepsilon_0 r).$$

Since $g_{A_i}^{L_i}(r)$ is monotone decreasing and integrable,

$$rg(r) \leq 2 \int_{r/2}^{r} g(s) \, ds \rightarrow 0$$

as $r \to \infty$, and we obtain $u_{\infty}(V') = u_{\infty}(V) \equiv u_{\infty}$. Accepting the following lemma, we proceed our proof of Theorem 1.2.

Lemma 2.1. Let
$$u_1(X) = u(X) - u_{\infty}$$
. Then
(2.12) $s - \lim_{t \to \pm \infty} (e^{iu_1(t,x)} - 1)e^{-itH_0}f = 0$

for all $f \in \mathcal{H}$.

Now we define a unitary propagator

$$\widetilde{U}(t,s):=e^{-iu_1(t,x)}U_{A(2)}(t,s)e^{iu_1(s,x)}.$$

Then, taking account of (2.10), we can easily verify that

$$i\frac{d}{dt}\widetilde{U}(t,s) = H_{A(1)}(t)\widetilde{U}(t,s), \quad \widetilde{U}(s,s) = I.$$

This means $\widetilde{U} = U_{A(1)}$ by the uniqueness of the propagator. Hence, in virtue of Lemma 2.1,

$$W_{A(1)}^{\pm}(s) = s - \lim_{t \to \pm \infty} e^{-\iota u_1(s,x)} U_{A(2)}(s,t) \{ (e^{\iota u_1(t,x)} - 1) + 1 \} e^{-\iota(t-s)H_0}$$

= $e^{-\iota u_1(s,x)} W_{A(2)}^{\pm}(s),$

Thus we have proved $S_{A(1)}(s) = S_{A(2)}(s)$.

Proof of Lemma 2.1. We use the same notation as in the proof of Proposition 1.1. It suffices to show (2.12) for a fixed $f \in G$. Since $\operatorname{supp} \widehat{f} \cap K_A = \emptyset$, we can find $\varepsilon_0 \in (0, 1/2)$ such that

(2.13)
$$|V^{\pm}(\xi)_{L_{\tau}}| \leq (1-2\varepsilon_0)|V^{\pm}(\xi)|,$$

for $\xi \in \text{supp } \hat{f}, j=1, \dots, N$. Define

$$\Gamma_{j} := \{X \in \mathbb{R}^{4}; |X_{L_{j}}| \leq \delta_{0}|X|\}$$

for $j=1,\dots,N$, where δ_0 is a fixed constant with $0 < \delta_0 \leq \varepsilon_0$. Then

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$$\begin{aligned} |\langle X, V^{\pm}(\xi) \rangle| &\leq |X_{L}| |V^{\pm}(\xi)_{L}| + |X_{L}| |V^{\pm}(\xi)_{L}| \\ &\leq (1 - \varepsilon_{0}) |X| |V^{\pm}(\xi)| \end{aligned}$$

if $\xi \in \operatorname{supp} \widehat{f}$ and $X = (t,x) \in \Gamma_j$. Thus,

$$|\nabla_{\xi}g^{\pm}(t,x,\xi)| = |X - tV^{\pm}(\xi)|$$

= $|X|^{2} + |t|^{2}|V^{\pm}(\xi)|^{2} - 2t\langle X, V^{\pm}(\xi) \rangle$
 $\geq \varepsilon_{0}(|X|^{2} + |t|^{2}|V^{\pm}(\xi)|^{2})$
 $\geq \varepsilon_{0}(|X|^{2} + |t|^{2}).$

Thus, due to the same argument as in the proof of (2.8), it follows that for any n there exists K_n such that

(2.14)
$$|(e^{-ttH_0}f)(x)| \leq K_n(1+|t|+|x|)^{-n}.$$

if $|t| \ge 1$ and $(t,x) \in \Gamma$, for some j. Now

$$(e^{iu_{1}(t,x)}-1) (e^{-itH_{6}}f) (x)$$

$$= (e^{iu_{1}(t,x)}-1) \prod_{j=1}^{N} (\chi ((t,x) \in \Gamma_{j}) + \chi ((t,x) \in \Gamma_{j})) (e^{-itH_{6}}f) (x)$$

$$= (e^{iu_{1}(t,x)}-1) \cdot \prod_{j=1}^{N} \chi ((t,x) \notin \Gamma_{j}) \cdot (e^{-itH_{6}}f) (x)$$

$$+ (e^{iu_{1}(t,x)}-1) \left[\sum_{j=1}^{N} \widetilde{\chi}_{j} (t,x) \chi ((t,x) \in \Gamma_{j}) \right] (e^{-itH_{6}}f) (x)$$

$$\equiv I_{1}(t,x) + I_{2}(t,x),$$

where $\widetilde{\chi}_{j}(t,x)$ are products of characteristic functions. In view of (2.14) we have

$$(2.15) s - \lim_{t \to \pm \infty} I_2(t, \cdot) = 0$$

Let $X \notin \bigcup_{j=1}^{N} L_{j}^{\perp}$. Then

$$u_{\infty} = \lim_{r \to \infty} u(rX)$$

=
$$\lim_{r \to \infty} \int_{0}^{1} \langle rX, A(\tau rX) \rangle d\tau$$

=
$$\int_{0}^{\infty} \langle X, A(\tau X) \rangle d\tau.$$

and so

$$u_1(X) = \int_{\infty}^1 \langle X, A(\tau X) \rangle d\tau = \sum_{j=1}^N \int_{\infty}^1 \langle X, A_j(\tau X) \rangle d\tau.$$

Thus, noting that $(t,x) \notin \bigcup_{j=1}^{N} L_{j}^{\perp}$ if $(t,x) \notin \bigcup_{j=1}^{N} \Gamma_{j}$, we estimate

$$\begin{aligned} |e^{iu_{1}(t,x)} - 1| &\leq |u_{1}(t,x)| \\ &\leq \sum_{j=1}^{N} \int_{1}^{\infty} |(t,x)| |A_{j}(s(t,x))| ds \\ &\leq \sum_{j=1}^{N} \int_{1}^{\infty} \delta_{0}^{-1} |(t,x)_{L_{j}}| g_{A_{j}}^{L_{j}}(s|(t,x)_{L_{j}}|) ds \\ &\leq \sum_{j=1}^{N} \int_{|(t,x)_{L_{j}}|}^{\infty} \delta_{0}^{-1} g_{A_{j}}^{L_{j}}(s) ds \end{aligned}$$

for $(t,x) \notin \bigcup_{j=1}^{N} \Gamma_{j}$. Hence, taking account of (2.9) and $g_{A_{j}}^{L_{j}} \in L^{1}((0,\infty))$, we obtain

$$s - \lim_{t \to \pm \infty} I_1(t, \cdot) = 0,$$

and finish the proof of Lemma 2.1.

Remark. Let ϕ be a real-valued function on R_t with

$$|\phi(t)| \leq K(1+|t|)^{-1-\varepsilon}$$
 uniformly on \mathbb{R}_t

for some K>0 and $\varepsilon>0$, and let $A=(\phi,0)$. Then we see that $F_A=F_0=0$ and

$$U_A(t,s) = e^{-\iota(\Phi(t) - \Phi(s))} e^{-\iota(t-s)H_0}$$
, where $\Phi(t) = \int_{-\infty}^t \phi(s) ds$.

Thus it follows that

$$W_A^{\pm}(s) = e^{\imath \int_{-\infty}^{\pm \infty} \phi(t) dt - \imath \Phi(s)}$$

and so

$$S_A(s) = e^{-i \int_{-\infty}^{\infty} \phi(t) dt}.$$

This means that $S_A(0) = S_0(0) = I$ if and only if $\int_{-\infty}^{\infty} \phi(t) dt \in 2\pi \mathbb{Z}$.

The above argument shows that Theorem 1.2 does not hold in general if the condition dim $L_1 \ge 2$ is dropped out.

We next prove Theorem 1.3. To do so, we have to prepare some notations.

For a subspace L we define a set $C(L) \subset S^2$ by

$$C(L) := \{ \omega \in S^2; \ \widetilde{\omega} \in L^\perp \},\$$

where $\widetilde{\omega} := (1, c\omega) \in \mathbb{R}^4$. We sometimes write C(V) = C(L) if L is the onedimensional subspace spanned by $V \neq 0 \in \mathbb{R}^4$. If $V = (v_0, v) \in \mathbb{R} \times \mathbb{R}^3$, HIROSHI T. ITO

$$C(V) = \{ \omega \in S^2; \langle v, \omega \rangle_{R^3} = -v_0/c \}.$$

Hence, C(V) is a circle on S^2 if $V \in D$, one point $\{-v_0v/c|v|^2\}$ on S^2 if $V \in \overline{D} \setminus D$ and empty if $V \notin \overline{D}$. Moreover, it is easy to see that

$$(2.16) C(V_1) = C(V_2) \neq \emptyset \Longleftrightarrow V_1 // V_2.$$

On the other hand,

$$C(L) = \bigcap_{j=1}^{N} C(V_j)$$

if $\{V_1, \dots, V_N\}$ is a basis of L. Therefore, the number of elements of C(L) is at most two if dimL=2. Moreover, it is at most one and zero if dimL=3 and dimL=4, respectively, since dim $L^{\perp}=1$ and dim $L^{\perp}=0$, respectively.

j=1

For
$$A = \sum_{j=1}^{N} A_j \in S$$
 with $A_j \in S(L_j)$ we define
 $C_A := \bigcup_{j=1}^{N} C(L_j)$

Then line integrals

(2.17)
$$K_{A}^{\omega}(\eta) := \int_{-\infty}^{\infty} \langle \widetilde{\omega}, A(t\widetilde{\omega} + \eta) \rangle_{R^{4}} dt, \quad \eta \in \mathbb{R}_{t} \times \mathbb{R}_{x}^{3},$$

is well-defined if $\omega \in S^2 \setminus C_A$. The matrix $\alpha \cdot \omega := \sum_{j=1}^3 \alpha_j \omega_j$ for $\omega = (\omega_1, \omega_2, \omega_3) \in S^2$ has eigenvalues 1 and -1 with multiplicity two, respectively. Thus, $P_{\pm}(\omega) := \frac{1}{2}(1 \pm \alpha \cdot \omega)$ is the eigenprojection associated with the eigenvalue ± 1 of $\alpha \cdot \omega$.

The proof of Theorem 1.3 is based on the following proposition.

Proposition 2.2. Let $A \in S$ and $s \in \mathbb{R}$. Then

(2.18)
$$w - \lim_{E \to +\infty} e^{-iE\omega \cdot x} P_{\pm} S_A(s) P_{\pm} e^{iE\omega \cdot x} = e^{iK_{\pm}^{*\omega}(s,x)} P_{\pm}(\omega)$$

if $\pm \omega \in S^2 \setminus C_A$. Here, $e^{\pm i E \omega \cdot x}$ and $e^{i K_A^{\pm \omega}(s,x)}$ are multiplication operators.

Remark 1. We can also show that

(2.19)
$$w - \lim_{E \to +\infty} e^{-iE\omega \cdot x} P_{\pm} S_A(s) P_{\mp} e^{iE\omega \cdot x} = 0$$

for each $s \in \mathbb{R}$ if both ω and $-\omega$ belong to $S^2 \setminus C_A$.

Remark 2. If the potential A is time-independent and short-range,

Proposition 2.2 has been already proved by Jung [J], and his proof is applicable to our case with a slight modification.

Lemma 2.3. Let $f \in \mathcal{H}$. (i) Suppose $\omega \in S^2 \setminus C_A$. Then there exists $E_0 > 0$ such that

(2.20)
$$\lim_{t \to \pm \infty} \| (W_A^{\pm}(0) - U_A(0, t) e^{-tH_0}) P_+ e^{iE\omega \cdot x} f \| = 0$$

uniformly in $E \ge E_0$.

(ii) Suppose $-\omega \in S^2 \setminus C_A$. Then there exists $E_0 > 0$ such that

(2.21)
$$\lim_{t \to \pm \infty} \| (W_A^{\pm}(0) - U_A(0,t) e^{-itH_0}) P_- e^{iE\omega \cdot x} f \| = 0$$

uniformly in $E \ge E_0$.

The proof is very similar to that of Proposition 1.1.

Proof. We use the notation in the proof of Proposition 1.1 and prove (i) for $f \in G$, since G is dense in \mathcal{H} . It suffices to show the following estimate (see (2.3)):

(2.22)
$$|(e^{-itH_0}P_+e^{iE\omega \cdot x}f)(x)| \le C_N(1+|t|+|x|)^{-N}$$

on Ω_j uniformly in $E \ge E_0$. We write

$$(e^{-\iota tH_0}P_+e^{\iota E\omega \cdot x}f)(x) = (2\pi)^{-3/2} \int e^{\iota g^+(t,x,\xi)}P^+(\xi) \ \hat{f}(\xi - E\omega) d\xi$$
$$= (2\pi)^{-3/2} \int e^{\iota g^+(t,x,\xi + E\omega)}P^+(\xi + E\omega) \ \hat{f}(\xi) d\xi.$$

Since $\widetilde{\omega}_{L_i} \neq 0$ for $j=1,\dots,N$ and $\lim_{E \to +\infty} V^+(\xi + E\omega) = \widetilde{\omega}$ uniformly in $\xi \in \operatorname{supp} \widehat{f}$, we can find $E_0 > 0$ and $\delta'_1 > 0$ such that

$$|V^+(\xi + E\omega)_{L_j}| \ge \delta'_1, \quad j=1,\cdots,N,$$

for any $E \ge E_0$ and $\xi \in \text{supp } \widehat{f}$. Moreover, we can verify that $|\partial_{\xi}^{\gamma}P^+(\xi + E\omega)|$ is uniformly bounded in $\xi \in \text{supp } \widehat{f}$ and in $E > E_0$ for each γ . Thus we can prove (2.22) in the same way as in the proof of (2.8) by using δ'_1 in place of δ_1 . \Box

Lemma 2.4. Let $f, g \in \mathcal{H}$ and let $\pm \omega \in S^2 \setminus C_A$. Then

$$\lim_{t \to +\infty} (e^{-iE\omega \cdot x} P_{\pm} e^{itH_0} U_A(t, -t) e^{itH_0} P_{\pm} e^{iE\omega \cdot x} f, g)$$
$$= (e^{-iE\omega \cdot x} P_{\pm} S_A(0) P_{\pm} e^{iE\omega \cdot x} f, g)$$

uniformly in $E \ge E_0$.

Proof. Denoting $W(t) := U_A(0,t) e^{-ttH_0}$, we have

$$\begin{aligned} |(e^{-iE\omega \cdot x}P_{\pm}W(t)^*W(-t)P_{\pm}e^{iE\omega \cdot x}f,g) \\ &-(e^{-iE\omega \cdot x}P_{\pm}S_A(0)P_{\pm}e^{iE\omega \cdot x}f,g)| \\ &\leq |(W(-t)P_{\pm}e^{iE\omega \cdot x}f,(W(t)-W_A^+(0))P_{\pm}e^{iE\omega \cdot x}g)| \\ &+|((W(-t)-W_A^-(0))P_{\pm}e^{iE\omega \cdot x}f,W_A^+(0)P_{\pm}e^{iE\omega \cdot x}g)|. \end{aligned}$$

Hence, using Lemma 2.3, we get the desired result. \Box

We define

$$Q(t,x) := H_A(t) - H_0 = -c \sum_{j=1}^{3} \alpha_j A^j(t,x) - A^0(t,x) I_4,$$
$$W_{\omega}(t,x) := \langle \widetilde{\omega}, A(t,x+ct\omega) \rangle_{R^4} P_+(\omega) + \langle (-\omega), A(t,x-ct\omega) \rangle_{R^4} P_-(\omega).$$

Lemma 2.5. For each
$$t > 0$$
 and each $\omega \in S^2$, one has

$$(2.23) \qquad s - \lim_{E \to +\infty} e^{-iE\omega \cdot x} e^{itH_0} U_A(t, -t) e^{itH_0} e^{iE\omega \cdot x} = e^{i\int_{-t}^{t} W_{\omega}(s, x) ds}$$

This lemma can be proved in the same way as in the proof of Theorem 3.1 in [J]. So, we only give a sketch of the proof.

Sketch of the proof. The Dyson expansion

(2.24)
$$e^{-iE\omega \cdot x} e^{itH_0} U_A(t, -t) e^{itH_0} e^{iE\omega \cdot x}$$
$$= I + \sum_{n=1}^{\infty} (-i)^n \int_{-t < t_1 < \cdots < t_n < t} dt_n \cdots dt_1 R_E(t_n) \cdots R_E(t_1)$$

converges in the operator norm of $\mathscr H$ uniformly in large E > 0, where

$$R_E(t):=e^{-tE\omega\cdot x}e^{ttH_0}Q(t,\cdot)e^{-ttH_0}e^{tE\omega\cdot x}.$$

Using

(2.25)
$$s - \lim_{E \to +\infty} \left[e^{-iE\omega \cdot x} e^{-isH_0} e^{iE\omega \cdot x} - e^{-isc(E+\omega \cdot D)\alpha \cdot \omega} \right] = 0,$$

we can see that

(2.26)
$$s - \lim_{E \to +\infty} \int_{-t}^{t} R_E(s) ds = - \int_{-t}^{t} W_\omega(s, \cdot) ds.$$

In the same way we also have

(2.27)
$$s - \lim_{E \to +\infty} \int_{-t < t_1 < \cdots < t_n < t} dt_n \cdots dt_1 R_E(t_n) \cdots R_E(t_1)$$

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$$= (-1)^n \int_{-t < t_1 < \cdots < t_n < t} dt_n \cdots dt_1 W_\omega(t_n, \cdot) \cdots W_\omega(t_1, \cdot)$$
$$= \frac{(-1)^n}{n!} \left(\int_{-t}^t W_\omega(s, \cdot) ds \right)^n,$$

where we have used $[W_{\omega}(t, \cdot), W_{\omega}(s, \cdot)] = 0$ for $s, t \in \mathbb{R}$. Thus (2.23) follows from (2.24) and (2.27).

In the momentum space it can be easily seen that

(2.28)
$$s - \lim_{E \to +\infty} e^{-iE\omega \cdot x} P_{\pm} e^{iE\omega \cdot x} = \frac{1}{2} (I \pm \alpha \cdot \omega) = P_{\pm}(\omega).$$

From this together with Lemma 2.5 it follows that

(2.29)
$$\lim_{E \to +\infty} e^{-iE\omega \cdot x} P_{\pm} e^{itH_0} U_A(t, -t) e^{itH_0} P_{\pm} e^{iE\omega \cdot x}$$
$$= e^{i \int_{-t}^{t} \langle (\pm \omega), A(s, x \pm cs\omega) \rangle_{B^*} ds} P_{\pm}(\omega).$$

Proof of Proposition 2.2. Let $f, g \in \mathcal{H}$ and define

$$S_E^{\pm}(t) := e^{-\iota E \omega \cdot x} P_{\pm} e^{\iota t H_0} U_A(t, -t) e^{\iota t H_0} P_{\pm} e^{\iota E \omega \cdot x}$$

and write

$$(e^{-iE\omega \cdot x}P_{\pm}S_{A}(0)P_{\pm}e^{iE\omega \cdot x}f,g) - (e^{iK_{h}^{\pm\omega}(0,x)}P_{\pm}(\omega)f,g)$$

= $((e^{-iE\omega \cdot x}P_{\pm}S_{A}(0)P_{\pm}e^{iE\omega \cdot x} - S_{E}^{\pm}(t))f,g)$
+ $((S_{E}^{\pm}(t) - e^{i\int_{-t}^{t}\langle\langle\langle \pm \omega\rangle\rangle,A(s,x\pm cs\omega)\rangle_{R}ds}P_{\pm}(\omega))f,g)$
+ $((e^{i\int_{-t}^{t}\langle\langle \pm \omega\rangle,A(s,x\pm cs\omega)\rangle_{R}ds} - e^{iK_{h}^{\pm\omega}(0,x)})P_{\pm}(\omega)f,g).$

The first term tends to zero as $t \to +\infty$ uniformly in large E > 0 by Lemma 2.4 and the second term tends to zero as $E \to +\infty$ for each t > 0 by (2.29). Since $\pm \omega \in S^2 \setminus C_A$, the function $A(s, x \pm csw)$ of s is integrable for each $x \in \mathbb{R}^3$. Therefore, the third term goes to zero as $t \to +\infty$. Hence, we have proved Proposition 2.2 for s = 0. Recalling (1.2) and noting that (see (2.25) and (2.28))

$$s - \lim_{E \to +\infty} \left[e^{-iE\omega \cdot x} P_{\pm} e^{isH_0} e^{iE\omega \cdot x} - P_{\pm}(\omega) e^{\pm is(cE + c\omega \cdot D)} \right] = 0,$$

we have

$$w - \lim_{E \to +\infty} e^{-iE\omega \cdot x} P_{\pm} S_A(s) P_{\pm} e^{iE\omega \cdot x}$$

= $P_{\pm}(\omega) e^{\pm isc\omega \cdot D} e^{iK_{\pm}^{\pm \omega}(0,x)} e^{\pm isc\omega \cdot D} P_{\pm}(\omega)$
= $e^{iK_{\pm}^{\pm \omega}(0,x\pm sc\omega)} P_{\pm}(\omega)$
= $e^{iK_{\pm}^{\pm \omega}(s,x)} P_{\pm}(\omega)$.

This completes the proof. \Box

Here we give the idea of the proof of Theorem 1.3 for a simple case, $A = (\phi, 0), \phi \in \mathcal{S}$. In this case, the right hand side of (2.18) determines

$$K^{\omega}_{A}(\eta) := \int_{-\infty}^{\infty} \phi(t\widetilde{\omega} + \eta) dt,$$

since $K_A^{\omega}(\eta) \to 0$ as $|\eta| \to \infty$ with $\langle \eta, \widetilde{\omega} \rangle = 0$. Thus the Fourier transform and the inverse Fourier transform yield

$$\widehat{\phi}(\varXi) = (2\pi)^{-2} \sqrt{1+c^2} \int_{\Pi_{\widetilde{\omega}}} e^{-i\langle\eta, \varXi\rangle} K_A^{\omega}(\eta) d\eta, \quad \text{for} \quad \varXi \in \Pi_{\widetilde{\omega}},$$
$$K_A^{\omega}(\eta) = (2\pi\sqrt{1+c^2})^{-1} \int_{\Pi_{\widetilde{\omega}}} e^{i\langle\eta, \varXi\rangle} \widehat{\phi}(\varXi) d\varXi, \quad \text{for} \quad \eta \in \Pi_{\widetilde{\omega}},$$

where $\Pi_{\theta} := \{\eta \in \mathbb{R}^4; \langle \theta, \eta \rangle_{\mathbb{R}^4} = 0\}$. On the other hand, we can easily verify that

$$\bigcup_{\omega \in S^2} \prod_{\widetilde{\omega}} = \overline{D}, \text{ the closure of } D.$$

Therefore it follows that the only $\phi |_{\overline{D}}$ is reconstructed from the right hand side of (2.18).

Remark. The potential ϕ can be completely reconstructed in the case of Schrödinger equations by a similar method [We]. Indeed, we obtain line integrals $\int_{-\infty}^{\infty} \phi(\eta_0, t\omega + \eta') dt$, $\eta = (\eta_0, \eta') \in \mathbb{R} \times \mathbb{R}^3$, instead of $K_A^{\omega}(\eta)$, from the scattering operator in the case of Schrödinger equations. Then, using the same argument as above, we see that these line integrals determine the whole $\hat{\phi}$. Hence the potential ϕ is completely reconstructed.

Now we define

$$\widetilde{C}_A := C_A \cup \Big(\bigcup_{k=1}^{N'} C(L'_k)\Big).$$

Lemma 2.6. Let $\omega \in S^2 \setminus \widetilde{C}_A$ and $\theta \in S^3$ such that $\widetilde{\omega}$ and θ are linearly independent, and let $\eta \in \mathbb{R}^4$. Then one has

(2.30)
$$\frac{d}{ds}K_{A}^{\omega}(s\theta+\eta)|_{s=0} = -\int_{-\infty}^{+\infty} \langle \widetilde{\omega} \wedge \theta, F_{A}(t\widetilde{\omega}+\eta) \rangle_{R^{*}} dt,$$

where $X \wedge Y = (X_1 Y_k - X_k Y_j)_{0 \le j < \kappa \le 3} \in \mathbb{R}^6$ for $X = (X_0, \dots, X_3)$ and $Y = (Y_0, \dots, Y_3)$.

Remark. Since $\widetilde{\omega} \in S^2 \setminus \widetilde{C}_A$, $F_A(t\widetilde{\omega} + \eta)$ is integrable with respect to $t \in \mathbb{R}$.

Proof. By Stokes' theorem we have

$$K_A^{\omega}(\eta) - K_A^{\omega}(\eta + s_0\theta) = \int_0^{s_0} ds \int_{-\infty}^{\infty} \langle \widetilde{\omega} \wedge \theta, F_A(t\widetilde{\omega} + s\theta + \eta) \rangle_{R^0} dt,$$

from which the lemma follows immediately.

Since

$$\frac{d}{ds}\exp(-iK_A^{\omega}(s\theta+\eta)) = -i\frac{d}{ds}K_A^{\omega}(s\theta+\eta)\cdot\exp(-iK_A^{\omega}(s\theta+\eta)),$$

we can conclude from Proposition 2.2 and Lemma 2.6 that for each $\eta \in \mathbb{R}^4$, $\omega \in S^2 \setminus \widetilde{C}_A$ and $\theta \in S^3$ the integral

(2.31)
$$\int_{-\infty}^{+\infty} \langle \widetilde{\omega} \wedge \theta, F_A(t\widetilde{\omega} + \eta) \rangle_{R^6} dt$$

can be constructed from the scattering operator $S_A(0)$ (see(1.2)).

Proof of Theorem 1.3. For a while we fix $E_0 = (\tau_0, \xi_0) \in D \setminus \Sigma$. Since $C(E_0) \cap \widetilde{C}_A$ is a finite or empty set due to (2.16), we can take $\{\omega_j^0\}_{j=1}^3 \subset C(E_0) \setminus \widetilde{C}_A$ with $\omega_j^0 \neq \omega_k^0$ if $j \neq k$. Then $\{\widetilde{\omega}_j^0\}_{j=1}^3$ are linearly independent in \mathbb{R}^4 , where $\widetilde{\omega}_j^0 = (1, c\omega_j^0)$. We also take $\widetilde{\omega}_0^0$ from \mathbb{R}^4 so that $\{\widetilde{\omega}_j^0\}_{j=0}^3$ is a basis of \mathbb{R}^4 . (Here we abused the notation $\widetilde{\omega}_0^0$, which need not be expressed as $(1, c\omega_0^0)$ for some $\omega_0^0 \in S^2$, to simplify notations below.) It should be notice that $\{\widetilde{\omega}_j^0 \wedge \widetilde{\omega}_k^0\}_{0 \leq k < j \leq 3}$ is also a basis of \mathbb{R}^6 .

We first assume $F_A \in L^1(\mathbb{R}_t \times \mathbb{R}^3_x)$. Noting that $\langle \Xi_0, \widetilde{\omega}_j^0 \rangle_{\mathbb{R}^4} = 0$ for j=1,2,3, we have, for each $0 \le k \le j \le 3$,

$$\begin{split} \langle \widetilde{\omega}_{j}^{0} \wedge \widetilde{\omega}_{k}^{0}, \widehat{F}_{A}(\Xi_{0}) \rangle_{R^{6}} \\ &= (2\pi)^{-2} \sqrt{1+c^{2}} \int_{\Pi \widetilde{\omega}_{j}^{0}} e^{-\iota \langle \eta, \Xi_{0} \rangle} d\eta \int_{-\infty}^{\infty} \langle \widetilde{\omega}_{j}^{0} \wedge \widetilde{\omega}_{k}^{0}, F_{A}(t \widetilde{\omega}_{j}^{0} + \eta) \rangle_{R^{6}} dt \end{split}$$

Since $\{\widetilde{\omega}_{j}^{0} \wedge \widetilde{\omega}_{k}^{0}\}_{0 \leq k < j \leq 3}$ is a basis of \mathbb{R}^{6} and since the integral with respect to t in the right hand side is determined by the scattering operator, $\widehat{F}_{A}(\Xi_{0})$ is determined by the scattering operator for each $\Xi_{0} \in D \setminus \Sigma$.

If F_A does not belong to $L^1(\mathbb{R}_t \times \mathbb{R}_x^3)$, we can not use the above formula directly since $\widehat{F_A}$ should be regarded as a distribution. So, more complicated arguments are needed to represent $\widehat{F_A}$ in terms of the line integrals (2.31) as in [St].

We set c=1 for simplicity, and fix j=1,2,3 for a while. Since ω_j^0 and ξ_0 are not parallel, they span a two-dimensional subspace, denoted by Λ_j , in \mathbf{R}_{ξ}^3 . We denote by Π_j the orthogonal complement of Λ_j in \mathbf{R}_{ξ}^3 , $\mathbf{R}_{\xi}^3=\Lambda_j\oplus\Pi_j$, and denote by ξ' , ξ'' the Λ_j , Π_j -component of $\xi \in \mathbf{R}_{\xi}^3$, respectively.

We can find a unique $\chi_j^0 \in \Lambda_j$, such that $\{\omega_j^0, \chi_j^0\}$ is an orthonormal basis of

 Λ_j with $\langle \chi_j^0, \xi_0 \rangle_{R^3} > 0$. Then we can write $\xi_0 = -\tau_0 \omega_j^0 + \lambda_j^0 \chi_j^0$, where $\lambda_j^0 := |\xi_0 + \tau_0 \omega_j^0| > 0$.

Now we define a family of orthonormal bases $\{\omega_j(\theta), \chi_j(\theta)\}, -\pi/4 < \theta < \pi/4$, of Λ_j by rotating the basis $\{\omega_j^0, \chi_j^0\}$:

$$\omega_{j}(\theta) := \omega_{j}^{9} \cos\theta + \chi_{j}^{9} \sin\theta,$$

$$\chi_{j}(\theta) := -\omega_{j}^{9} \sin\theta + \chi_{j}^{9} \cos\theta$$

It is easily seen that if (τ,ξ) is sufficiently near (τ_0,ξ_0) there exists a unique $\theta_j(\tau,\xi) \in \mathbb{R}$ near 0 satisfying $\langle \xi', \omega_j(\theta_j(\tau,\xi)) \rangle_{R^3} = -\tau$. We set $\lambda_j(\tau,\xi)$ $:= \langle \xi', \chi_j(\theta_j(\tau,\xi)) \rangle_{R^3} = \sqrt{|\xi'|^2 - \tau^2}$ and $\rho_j(\tau,\xi) := \langle \xi'', \overline{\omega}_j^0 \rangle$, where $\overline{\omega}_j^0$ is a unit vector in Π_j .

Consequently we can see that there are a small neighborhood U_0 of $(\tau_0, \xi_0) \in \mathbb{R} \times \mathbb{R}^3$, a small neighborhood $V_j \subset \{(\theta_j, \lambda_j, \rho_j, \tau_j) \in (-\pi/4, \pi/4) \times \mathbb{R} \times \mathbb{R} \times \mathbb{R}\}$ of $(0, \lambda_j^0, 0, \tau_0)$ and a diffeomorphism

$$(\theta_j, \lambda_j, \rho_j, \tau_j): U_0 \longrightarrow V_j$$

such that

$$\begin{aligned} \theta_j(\tau_0,\xi_0) &= 0, \quad \lambda_j(\tau_0,\xi_0) = \lambda_j^0, \quad \rho_j(\tau_0,\xi_0) = 0, \quad \tau_j(\tau_0,\xi_0) = \tau_0, \\ \tau_j(\tau,\xi) &= \tau, \quad -\tau_j(\tau,\xi) \, \omega_j(\theta_j(\tau,\xi)) + \lambda_j(\tau,\xi) \, \chi_j(\theta_j(\tau,\xi)) + \rho_j(\tau,\xi) \, \varpi_j^0 = \xi. \end{aligned}$$

We also note that $|\partial(\tau,\xi)/\partial(\theta_j,\lambda_j,\rho_j,\tau_j)| = \lambda_j$.

Since $\{\widetilde{\omega}_k^0 \land \widetilde{\omega}_j^0\}_{0 \le k < j \le 3}$ is a basis of \mathbb{R}^6 and since each function $\omega_j(\theta_j(\tau,\xi))$ is smooth on U_0 with $\omega_j(\theta_j(\tau_0,\xi_0)) = \omega_j^0, \{\widetilde{\omega}_k^0 \land \widetilde{\omega}_j(\theta_j(\tau,\xi))\}_{0 \le k < j \le 3}$ is also a basis of \mathbb{R}^6 for each (τ,ξ) in a small neighborhood of (τ_0,ξ_0) , which can be assumed to be U_0 , where $\widetilde{\omega}_j(\theta) := (1,\omega_j(\theta))$. We can also assume that $\omega_j(\theta_j(\tau,\xi)) \notin \widetilde{C}_A$ for $(\tau,\xi) \in U_0$, $1 \le j \le 3$. Thus, for any $\phi \in C_0^\infty(U_0;\mathbb{C}^6)$ there exist $\phi_{jk}(\tau,\xi) \in C_0^\infty(U_0;\mathbb{C}), \ 0 \le k < j \le 3$, such that

$$\phi(\tau,\xi) = \sum_{0 \le k < j \le 3} \phi_{jk}(\tau,\xi) \,\widetilde{\omega}_k^0 \wedge \widetilde{\omega}_j(\theta_j(\tau,\xi))$$

on U_0 . So we have

$$\begin{split} &\langle \widehat{F}_{A}, \phi \rangle = \langle F_{A}, \widehat{\phi} \rangle \\ &= \iint dt dx \langle F_{A}(t,x), (2\pi)^{-2} \iint e^{-it\tau - i\langle x, \xi \rangle} \phi(\tau, \xi) d\tau d\xi \rangle_{C^{6}} \\ &= (2\pi)^{-2} \sum_{0 \le k < j \le 3} \iint dt dx \langle F_{A}(t,x), \iint e^{-it\tau - i\langle x, \xi \rangle} \phi_{jk}(\tau, \xi) \widetilde{\omega}_{k}^{0} \wedge \widetilde{\omega}_{j}(\theta_{j}(\tau, \xi)) d\tau d\xi \rangle_{C^{6}} \\ &= (2\pi)^{-2} \sum_{0 \le k < j \le 3} \iint dt dx \langle F_{A}(t,x), \langle f_{A}(t,x), f_{A}($$

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$$\iiint e^{-it\tau_j - i\langle x, \xi(\theta_j, \lambda_j, \rho_j, \tau_j) \rangle} \phi_{jk}(\tau_j, \xi(\theta_j, \Lambda_j, \rho_j, \tau_j)) \widetilde{\omega}_k^0 \wedge \widetilde{\omega}_j(\theta_j) \lambda_j d\tau_j d\rho_j d\lambda_j d\theta_j \rangle_{C^0}$$
$$= (2\pi)^{-2} \sum_{0 \le k < j \le 3} \iint dt dx \int G_{jk}(t, x, \theta) d\theta_j,$$

where

$$(2.32) \qquad \begin{array}{l} G_{jk}^{0}(t,x,\theta_{j}) := \langle F_{A}(t,x), \widetilde{\omega}_{k}^{0} \wedge \widetilde{\omega}_{j}(\theta_{j}) \rangle_{C^{6}}, \\ (2.32) \qquad G_{jk}^{1}(t,x,\theta_{j}) := \iiint e^{-\iota t\tau_{j} - \iota \langle x, \xi(\theta_{j}, \lambda_{j}, \theta_{j}, \tau_{j}) \rangle} \phi_{jk}(\tau_{j}, \xi(\theta_{j}, \lambda_{j}, \rho_{j}, \tau_{j})) \lambda_{j} d\tau_{j} d\rho_{j} d\lambda_{j}, \\ G_{jk}(t,x,\theta_{j}) = G_{jk}^{0}(t,x,\theta_{j}) G_{jk}^{1}(t,x,\theta_{j}). \end{array}$$

Integrating by parts with respect to $(\tau_i, \rho_i, \lambda_i)$ in (2.32), we can see that for any n > 0 there exists K_n such that

$$|G_{jk}(t,x,\theta_j)| \leq K_n(1+|t-\langle \omega_j(\theta_j),x\rangle|+|\langle \chi_j(\theta_j),x\rangle|+|\langle \varpi_j^0,x\rangle|)^{-n}|F_A(t,x)|$$

uniformly in (t, x, θ_j) . We will show that.

(2.33)
$$\sup_{\theta_j} \|G_{jk}(\cdot,\cdot,\theta_j)\|_{L^1(R_t\times R^3_x)} < +\infty.$$

Let $X(\theta_j)$ be the orthogonal projection of X = (t,x) onto $\prod_{\tilde{\omega}_j(\theta_j)}$. Then, observing that

$$\{(1, -\omega_j(\theta_j)), (0, \chi_j(\theta_j)), (0, \overline{\omega}_j^0)\}$$

is a basis of $\prod_{\widetilde{\omega}_{j}(\theta_{j})}$, we have

(2.34)
$$|G_{jk}(t,x,\theta_{j})| \leq K_{n}'(1+|X(\theta_{j})|)^{-n}|F_{A}(X)|$$
$$\leq K_{n}'(1+|X(\theta_{j})|)^{-n}\sum_{k=1}^{N'}g_{F_{k}}^{L'_{k}}(|X_{L'_{k}}|)$$

uniformly in (t, x, θ_j) for some $K'_n > 0$. Here note that $L'_k \subset \prod_{\tilde{\omega}, (\theta_j)}$ for all k since $\omega_j (\theta_j (\tau, \xi)) \notin \widetilde{C}_A$ if $(\tau, \xi) \in U_0$, $1 \le j \le 3$. Thus, (2.34) implies (2.33). Thus, Fubini's theorem guarantees that

(2.35)
$$\iint dt dx \int G_{jk}(t,x,\theta_j) d\theta_j = \int d\theta_j \iint G_{jk}(t,x,\theta_j) dt dx.$$

Now for each θ_j we introduce new coordinates $\{s_j, \lambda'_j, \rho'_j, \tau'_j\}$ for $\{t, x\}$:

 $t = \tau'_j + s_j, \qquad x = \tau'_j \omega_j(\theta_j) + \lambda'_j \chi_j(\theta_j) + \rho'_j \overline{\omega}_j^0$

where $s_j, \tau'_j, \lambda'_j, \rho'_j \in \mathbf{R}$. Then

$$e^{-\imath t\tau - \imath x \cdot \xi(\theta_i, \lambda_j, \rho_j, \tau_j)} = e^{-\imath (s_i \tau + \lambda'_j \lambda_j + \rho_j \rho_j)}$$

is independent of τ'_{j} , and so is G^1_{jk} . Thus, the right hand side of (2.35) is written as

$$\int d\theta_j \iiint G_{jk}^1 ds_j d\lambda'_j d\rho'_j \int_{-\infty}^{\infty} \langle F_A(\tau'_j + s_j, \tau'_j \omega_j(\theta_j) + \lambda'_j \chi_j(\theta_j) + \rho'_j \overline{\omega}_j^0), \widetilde{\omega}_k^0 \wedge \widetilde{\omega}_j(\theta_j) \rangle_{C^6} d\tau'_j.$$

Therefore we have

$$\begin{split} &\langle \widehat{F}_{A}, \phi \rangle \\ &= (2\pi)^{-2} \sum_{0 \le k < j \le 3} \int d\theta_{j} \iiint G_{jk}^{1} ds_{j} d\lambda'_{j} d\rho'_{j} \\ &\times \int_{-\infty}^{\infty} \langle F_{A}(\tau'_{j} \widetilde{\omega}_{j}(\theta_{j}) + \eta(\lambda'_{j}, s_{j}, \rho'_{j}, \theta_{j})), \widetilde{\omega}_{k}^{0} \wedge \widetilde{\omega}_{j}(\theta_{j}) \rangle_{C^{0}} d\tau'_{j}, \end{split}$$

where $\eta(\lambda'_i, s_i, \rho'_i, \theta_i) = (s_i, \lambda'_i \chi_i(\theta_i) + \rho'_i \overline{\omega}_i^0)$. Since $(\tau_0, \xi_0) \in D \setminus \Sigma$ is arbitrary and G^1_{jk} is independent of F_A , we have completed the proof of Theorem 1.3. \Box

§ 3. Proofs of Theorems 1.5, 1.7 and Corollary 1.6

Lemma 3.1. Suppose that F is a bounded function on \mathbb{R}^d , $d \ge 2$, such that

$$|F(x)| \le g(|\langle V, x \rangle_{R^d}|) \quad on \quad \mathbb{R}^d$$

for some $V \in S^{d-1}$ and g with $g(t)(1+t)^{-1} \in L^1((0,\infty))$. Let $\chi_{\varepsilon}(\xi) := \chi(\xi/\varepsilon), \varepsilon > 0$, for $\chi \in \mathscr{S}(\mathbb{R}^d)$ and let \widehat{F} be the Fourier transform of F. Then

(3.2)
$$\lim_{\varepsilon \downarrow 0} \langle \widehat{F}, \chi_{\varepsilon} u \rangle = 0$$

for any $u \in C_0^{\infty}(\mathbb{R}^d)$.

Proof. We assume V = (1,0,0,0) for simplicity and calculate

(3.3)

$$\langle \widehat{F}, \chi_{\varepsilon} u \rangle = \langle F, \widehat{\chi_{\varepsilon} u} \rangle$$

$$= \text{const.} \langle F, \widehat{\chi_{\varepsilon}} * \widehat{u} \rangle$$

$$= \text{const.} \iint \varepsilon^{d} F(x) \, \widehat{\chi} \, (\varepsilon(x-y)) \, \widehat{u} \, (y) \, dy \, dx.$$

Using the estimate $|\widehat{\chi}(x)| \leq K(1+|x|)^{-d}$, we get $|\varepsilon^d \widehat{\chi}(\varepsilon x)| \leq K(1+|x|)^{-d}$ uniformly in $\varepsilon \in (0,1)$. Therefore it follows that

$$|\varepsilon^{d}F(x)\,\widehat{\chi}\,(\varepsilon(x-y))\,\widehat{u}\,(y)|$$

$$\leq K|F(x)|(1+|x-y|)^{-d}|\,\widehat{u}\,(y)|$$

$$\leq Kg(|x_{l}|)\,(1+|x|)^{-d}(1+|y|)^{d}|\,\widehat{u}\,(y)|,$$

where we have used (3.1) in the last step. On the other hand, denoting $x = (x_1, x')$, we see that

$$\int (1+|x_1|+|x'|)^{-d} dx' \leq \text{const.} (1+|x_1|)^{-1}.$$

Thus, $g(|x_1|)(1+|x|)^{-d}(1+|y|)^d | \hat{u}(y)| \in L^1(\mathbb{R}^d_x \times \mathbb{R}^d_y)$. By applying Lebesgue's convergence theorem to (3.3), (3.2) follows.

Proof of Theorem 1.5. Let $\phi \in C_0^{\infty}(\mathbb{R}^4 \setminus \{0\}; \mathbb{C}^6)$. Due to (1.5), we can take $\chi \in C^{\infty}(\mathbb{R}^4)$ such that

$$\operatorname{supp}\chi\phi\subset D$$
, $\operatorname{supp}\{(1-\chi)\phi\}\cap\operatorname{supp}\widehat{F_A}=\emptyset$.

Therefore,

$$\langle \widehat{F}_A, \phi \rangle = \langle \widehat{F}_A, \chi \phi \rangle.$$

Thus Theorem 1.3 guarantees that $\widehat{F_A}|_{R^{1}\setminus\{0\}}$ can be reconstructed from the scattering operator. Hence, by Lemma 3.1 we see that $\widehat{F_A}$ and so F_A can be reconstructed from scattering operator, since each F_k satisfies (3.1).

Proof of Corollary 1.6. If F_k satisfies (1.6), the support of $\widehat{F_K}$ is contained in the three-dimensional subspace

$$V_k^{\perp} := \{ \Xi \in \mathbb{R}^4 ; \langle V_k, \Xi \rangle = 0 \}$$

which satisfies $V_k^{\perp} \setminus \{0\} \subset D$ if $V_k \in T$.

Proof of Theorem 1.7. (1) Let L be the one-dimensional subspace spanned by V and fix $\Xi_0 \notin L$. Then it is easy to see that there exist a neighborhood U_0 of $\Xi_0, \sigma_0 \subset \mathbb{R}$ and $\varepsilon_0 > 0$ such that

$$U_0 \cap L = \emptyset$$
, $U_0 + \sigma V \subset D \setminus \Sigma$ if $|\sigma - \sigma_0| < \varepsilon_0$.

We define $\phi_{\sigma}(\Xi) := \phi(\Xi - \sigma V)$ for any $\phi \in C_0^{\infty}(U_0)$ and $\sigma \in \mathbb{R}$. Then $\phi_{\sigma} \in C_0^{\infty}(U_0 + \sigma V) \in C_0^{\infty}(D \setminus \Sigma)$ if $|\sigma - \sigma_0| < \varepsilon_0$. Therefore Theorem 1.3 says that $\langle \widehat{F}_A, \phi_{\sigma} \rangle$ can be reconstructed from $S_A(0)$. On the other hand,

$$f(\sigma) := \langle \widehat{F_A}, \phi_\sigma \rangle = \langle F_A e^{-\iota \langle X, \sigma V \rangle}, \ \widehat{\phi} \rangle$$

is an analytic function on $\{\sigma \in C; \operatorname{Re}\sigma \in \mathbb{R}, |\operatorname{Im}\sigma| < \delta\}$ by (1.8). Hence, $f(\sigma)$ is determined by $S_A(0)$ for any $\sigma \in \mathbb{R}$, and so is $f(0) = \langle \widehat{F}_A, \phi \rangle$, in particular. Since \mathcal{E}_0 is arbitrary, this implies that $\widehat{F}_A|_{R^{A} \setminus L}$ is determined by $S_A(0)$.

(2) Let L be the one-dimensional subspace spanned by V and take $\phi \in C^{\infty}(L^{\perp})$ with $\phi(\eta) = 1$ for $|\eta| > 2$ and $\phi(\eta) = 0$ for $|\eta| < 1$, and define $\phi_{\varepsilon}(\eta) = \phi(\eta/\varepsilon)$ for $\varepsilon > 0$. It suffices to show that

$$\lim_{\varepsilon \downarrow 0} \phi_{\varepsilon}(\eta) \widehat{F}_{A}^{jk} = \widehat{F}_{A}^{jk} \text{ in } \mathcal{D}'(\mathbf{R}^{4}), \quad 0 \le j \le k \le 3,$$

since $\operatorname{supp}(\phi_{\varepsilon} u) \subset \mathbb{R}^4 \setminus L$ for all $u \in C_0^{\infty}$. To do so we assume V = (1,0,0,0) for simplicity and write X = (t,x). In the same way as in the proof of Lemma 3.1 it suffices to show that

$$\iint e^{-\delta|t|} g(|\langle V', X\rangle|) (1+|x|)^{-3} dt dx$$

$$\leq \text{const.} \iint e^{-\delta|t|/2} g(|\langle V', X\rangle|) (1+|X|)^{-3} dt dx$$

is finite. But, this follows immediately. 🗌

§ 4. Proof of Proposition 1.4

Proof of Proposition 1.4. Since $\phi \in \mathcal{S}$, $S_{\lambda A}(0)$ is real-analytic in λ (cf.[Th], p.249]):

$$S_{\lambda A}(0) = \sum_{n=0}^{\infty} (-\lambda)^n S_n,$$

where each S_n is a bounded operator on \mathcal{H} and

$$S_0=I, \quad S_1=\int_{-\infty}^{\infty}e^{ttH_0}\phi(t,\cdot)e^{-ttH_0}dt,\cdots.$$

So, it suffices to show that $S_1 \neq 0$.

We recall that the 4×4 matrix $h_0(\xi) = \alpha \cdot \xi + \beta$, the symbol of H_0 , has the eigenvalues $\pm \sqrt{|\xi|^2 + 1}$ with multiplicity two, respectively, and the associated eigenprojections are denoted by $P^{\pm}(\xi)$, and that P_{\pm} denotes the orthogonal projection in \mathcal{H} onto the positive/negative energy subspace of H_0 : $P_{\pm} = \frac{1}{2} \left(I \pm \frac{H_0}{1 + 1} \right)$

$$\frac{H_0}{|H_0|}$$

Then we have

$$\mathcal{F}_{x \to \xi} P_{-} S_{1} P_{+} \mathcal{F}_{x \to \xi}^{-1} = \int_{-\infty}^{\infty} P^{-}(\xi) e^{-\iota t \sqrt{|\xi|^{2} + 1}} \widetilde{\phi}(t) e^{-\iota t \sqrt{|\xi|^{2} + 1}} P^{+}(\xi) dt$$

where $\widetilde{\phi}(t) = \mathscr{F}_{x \to \xi} \phi(t, \cdot) \mathscr{F}_{x \to \xi}^{-1}$. Therefore $\mathscr{F}_{x \to \xi} P_- S_1 P_+ \mathscr{F}_{x \to \xi}^{-1}$ is the integral operator with kernel

$$T(p,q) := (2\pi)^{-3/2} \int_{-\infty}^{\infty} e^{-it\sqrt{|p|^2+1}} (\mathcal{F}_{x \to \xi}\phi) (t,p-q) e^{-it\sqrt{|q|^2+1}} P^{-}(p) P^{+}(q) dt$$
$$= (2\pi)^{-1} \widehat{\phi} (\sqrt{|p|^2+1} + \sqrt{|q|^2+1}, p-q) P^{-}(p) P^{+}(q).$$

Thus the proposition follows from Lemma 4.1 below since $\widehat{\phi}(\tau,\xi) \neq 0$ for some $(\tau,\xi) \in D_I$.

Lemma 4.1. For each $(\tau, \xi) \in D_1$, $\xi \neq 0$, there exist p and q in \mathbb{R}^3 such that (i) $\xi = p - q$, (ii) $\tau = \sqrt{|p|^2 + 1} + \sqrt{|q|^2 + 1}$, (iii) $P^-(p)P^+(q) \neq 0$.

Proof of Lemma 3.1. We can write $2\sqrt{1+r^2} = \tau$ for some constant r > 0because of $\tau > \sqrt{|\xi|^2 + 4}$. Since $|\xi|/r < 2$, we can take ω , $\theta \in S^2$ with $\omega - \theta = \xi/r$, and we define $p = r\omega$, $q = r\theta$. Then it is clear that they satisfy (i) and (ii). We next show that they satisfy (iii). Suppose that $P^-(p)P^+(q) = 0$. Then $P^+(p) \supset$ $P^+(q)$. On the other hand, both operators have the same rank, rank $P^+(p) =$ rank $P^+(q) = 2$, and so they coincide: $P^+(p) = P^+(q)$. Hence noting that $P^+(u) =$ $(1/2)(I + h_0(u)/|h_0(u)|)$ for each $u \in \mathbb{R}^3$ and that $|h_0(p)| = |h_0(q)| = \sqrt{1+r^2}$, we obtain $\alpha \cdot (p-q) = 0$, and it follows that p = q. This contradicts with $\omega - \theta = \xi/r$ $\neq 0$. We have thus proved the lemma. \Box

§ 5. Relativistic Invariance

A Lorentz transformation $\Lambda: \mathbb{R}^4 \longrightarrow \mathbb{R}^4$ is a linear map preserving the Lorentz metric,

$$\langle X, X' \rangle_{LM} := c^2 x_0 x_0' - \sum_{j=1}^3 x_j x_j',$$

where $X = (x_0, \dots, x_3)$, etc. This condition is equivalent to

(5.1)
$${}^{t}AGA = G$$
, where $G := \begin{pmatrix} c^{2}I_{1} & O \\ O & -I_{3} \end{pmatrix}$.

A Poincaré transformation,

is the combination of a Lorentz transformation Λ and a space-time translation by $a \in \mathbb{R}^4$ and is associated with changing an inertial frame with coordinate $X = {}^t(t,x)$ into another inertial frame with coordinate $X' = {}^t(t',x')$ defined by (5.2). By this change the Dirac equation

$$i\frac{\partial}{\partial t}\Psi(t,x) = \left[c\sum_{j=1}^{3}\alpha_{j}\left(\frac{1}{i}\frac{\partial}{\partial x_{j}}-A^{j}(t,x)\right) + \alpha_{4}mc^{2}-A^{0}(t,x)I_{4}\right]\Psi(t,x)$$

in the old inertial frame is converted into the Dirac equation

$$i\frac{\partial}{\partial t'}\Psi'(t',x') = \left[c\sum_{j=1}^{3}\alpha_{j}\left(\frac{1}{i}\frac{\partial}{\partial x_{j}'}-A^{*j}(t',x')\right)+\alpha_{4}mc^{2}-A^{*0}(t',x')I_{4}\right]\Psi'(t',x')$$

in the new inertial frame. Each component of the potential $A^*(X')$ in the new frame is written as a linear combination of those of $A(\Lambda^{-1}(X'-a))$ (see, e.g., [Tha]). Let V_1, \dots, V_n be a basis of a subspace L. Then there exists a constant K>0 such that

$$K^{-1}|X_L| \leq \sum_{j=1}^n |\langle V_j, X \rangle| \leq K|X_L|.$$

Hence we can see that $A^* \in S({}^tA^{-1}L)$ if $A \in S(L)$. Each component $F_{A^*}^{*im}(X')$ of the electromagnetic field $F_{A^*}^{*}(X')$ in the new frame is also written as a linear combination of components $F_A^{jk}(\Lambda^{-1}(X'-a))$, $0 \le j \le k \le 3$, of $F_A(\Lambda^{-1}(X'-a))$. Thus, carrying out the Fourier transform, we have

$$\widehat{F_{A^*}^{*im}}(\varXi) = \sum_{0 \le j < k \le 3} c_{jk} e^{-\iota \langle \varXi, a \rangle} \widehat{F_A^{*i}}({}^t \Lambda^{-1} \varXi), \quad \varXi \in \mathbb{R}^4,$$

where c_{jk} 's are constants determined by Λ . Hence, $F_{\Lambda}^* \in \widetilde{S}({}^t\Lambda^{-1}L)$ if $F_A \in \widetilde{S}(L)$. On the other hand, by virtue of (5.1) we see that $\Lambda^{-1}G^{-1t}\Lambda^{-1} = G^{-1}$. Then it follows that

$$\langle {}^{t}\Lambda^{-1}\mathcal{E}, {}^{t}\Lambda^{-1}\mathcal{E}' \rangle_{LM^{*}} = \langle \mathcal{E}, \mathcal{E}' \rangle_{LM^{*}},$$

where

$$\langle E, E' \rangle_{LM^*} := \xi_0 \xi_0' - c^2 \sum_{j=1}^3 \xi_j \xi_j'$$

for $\Xi = (\xi_0, \dots, \xi_3)$, etc. Thus, we can see that

 ${}^{t}\Lambda^{-1}D = D.$

Namely, the set D in the dual space of the space-time is invariant under Poincaré transformations in the space-time. We also see that D_I in Proposition 1.4 is invariant.

After all we can conclude Theorems 1.2, 1.3, 1.5, 1.7 and Corollary 1.6 hold on any inertial frame.

In the last of this section we show that each field F_k in Corollary 1.6 is regarded as a time-independent one on a suitable inertial frame. Let $V = V_k$ is timelike;

$$V \in T := \{(t, x) \in \mathbb{R}^4; \quad c|t| > |x|\}$$

with $\langle V, V \rangle_{LM} = c^2$. Then it is known that there exists a Lorentz transform Λ : $\mathbb{R}_t \times \mathbb{R}^3_x \longrightarrow \mathbb{R}_{t'} \times \mathbb{R}^3_{x'}$ such that $\Lambda V = {}^t(1,0,0,0)$. Therefore $F_A(\Lambda^{-1} \cdot)$ is t'-independent.

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