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ESTIMATES OF ONE-DIMENSIONAL OSCILLATORY INTEGRALS

by Detlef MÜLLER

1. Introduction.

If U is an open domain in \mathbb{R}^k and if f is a smooth, real valued function on U, one may define the associated oscillatory integral as

$$E_f(\vartheta) = \int_{U} \vartheta(x) e^{2\pi i f(x)} dx,$$

where ϑ belongs to $\mathscr{D}(U)$, the space of testfunctions on U.

When f has the form $f = \sum_{j=1}^{n} \eta_{j} \psi_{j}$, where the $\psi_{j} \in C^{\infty}(U)$ are real-valued functions and η_{j} are real parameters, one is interested in the asymptotic behaviour of $E_{\Sigma \eta_{j} \psi_{j}}(\vartheta)$ as $(\eta_{1}, \ldots, \eta_{n})$ tends to infinity, for several reasons.

For example, if μ is a smooth measure on a smooth submanifold of \mathbf{R}^m , and if the support of μ is sufficiently small, then the Fourier-Stieltjes transform $\hat{\mu}(\eta_1,\ldots,\eta_n)$ may always be written as $E_{\Sigma\eta_j\psi_j}(9)$ for certain functions ψ_j and 9.

Good information about the asymptotic behaviour of such Fourier-Stieltjes transforms is needed to solve the synthesis problem for smooth submanifolds of \mathbf{R}^m (see e.g. [7]). And, as Professor Y. Domar has pointed out to me, such knowledge would also yield information about the decay at infinity of solutions of partial differential equations (see e.g. [5]).

As far as I know, satisfactory aswers to the above problem have only been given for oscillatory integrals $E_{\Sigma\eta,\psi_i}(\vartheta)$ with

$$\Sigma \eta_{j} \psi_{j}(x_{1}, \ldots, x_{k}) = \sum_{j=1}^{k} \eta_{j} x_{j} + \eta_{k+1} \psi_{k+1}(x_{1}, \ldots, x_{k}),$$

which correspond to surface carried measures (see [2], [4], [6]). In some sense, the other extreme is the case where $\Sigma \eta_j \psi_j$ is a function of only one real variable, which corresponds to measures on curves. For this case, we will prove some quite general results.

2.

Let $\psi \in C^{\infty}(I, \mathbf{R}^n)$, $\psi = (\psi_1, \dots, \psi_n)$, where $I \neq \emptyset$ is some bounded open interval in \mathbf{R} . For ξ , $\eta \in \mathbf{R}^n$ let $\xi \cdot \eta$ denote the Euclidean inner product on \mathbf{R}^n , and correspondingly let

$$\eta \cdot \psi(x) = \sum_{j=1}^{n} \eta_{j} \psi_{j}(x).$$

Further let

$$|\eta| := \max_{j} |\eta_{j}| \quad \text{for} \quad \eta \in \mathbf{R}^{n}.$$

Define the torsion τ of ψ by

$$\tau(x) = \det (\psi_j^{(i+1)}(x))_{i,j=1,...,n} = \det (\psi''(x)\psi'''(x)...\psi^{(n+1)}(x)),$$

where ψ is regarded as a column vector and $\psi^{(k)}$ denotes the k – th derivative of ψ . At least for n=2 we have $\tau(x)=k(x)|\psi''(x)|^2$, where k is the torsion of the curve $\gamma=\{(x,\psi(x)):x\in I\}$ in \mathbb{R}^{n+1} . Let

$$e(t) = e^{2\pi i t}$$
 for $t \in \mathbb{R}$, and $e(g) = e \circ g$

for $g \in C^{\infty}(I, \mathbb{R})$. If $\psi_0(x) = x$ for $x \in \mathbb{R}$, then for $\vartheta \in \mathcal{D}(I)$, $\eta_0 \in \mathbb{R}$ and $\eta = (\eta_1, \dots, \eta_n) \in \mathbb{R}^n$, we have

$$E_{n} \sum_{j=0}^{n} \eta_{j} \psi_{j} (\vartheta) = (\vartheta e(\eta \cdot \psi)) (-\eta_{0}).$$

So it will be slightly more general to study the behaviour of $|\vartheta e(\eta \cdot \psi)|_{PM}$ as $|\eta| \to \infty$, where

$$|\varphi|_{PM} = \sup_{t \in \mathbb{R}} |\hat{\varphi}(t)|$$

for every $\varphi \in \mathcal{D}(\mathbf{R})$.

For certain reasons (see [3]; [7], Th. 4.1), we will also study $|\vartheta e(\eta \cdot \psi)|_A$, where

$$|\varphi|_{A} = \int |\hat{\varphi}(t)| dt$$

for every $\varphi \in \mathcal{D}(\mathbf{R})$.

We will first state our main results and prove some corollaries:

Theorem 1. – Let $\vartheta \in \mathcal{D}(I)$. Then

- (i) $|\vartheta e(\eta \cdot \psi)|_A = 0(|\eta|^{\frac{1}{2}})$, as $|\eta| \to \infty$.
- (ii) If for some subinterval J of I and some $\sigma > 0$

$$|\vartheta(x)| \geqslant \sigma$$
 and $|\vartheta(x) - \vartheta(y)| < \sigma/2$ for all $x, y \in J$,

and if $\psi_1|_1,\ldots,\psi_n|_J$ are linearly independent modulo affine linear functions, then there is a constant C>0, such that

$$|\vartheta e(\eta \cdot \psi)|_{A} \geqslant C(1+|\eta|)^{\frac{1}{2}}$$

for all $\eta \in \mathbf{R}^n$.

COROLLARY 1. – The following two conditions are equivalent:

(i) For each $\vartheta \in \mathcal{D}(\mathbf{R})$, $\vartheta \neq 0$, there are constants c > 0, C > 0, such that for all $\eta \in \mathbf{R}^n$

$$c(1+|\eta|)^{\frac{1}{2}} \le |\vartheta e(\eta \cdot \psi)|_A \le C(1+|\eta|)^{\frac{1}{2}}.$$

(ii) ψ_1, \ldots, ψ_n are linearly independent modulo affine linear functions on every non empty open subinterval of I.

Proof of Corollary 1. — (i) follows directly from (ii) by Theorem 1. Now suppose that there exists a vector $v \in \mathbf{R}^n$, $v \neq 0$, such that $v \cdot \psi$ is affine linear on some open subinterval $\mathcal{J} \neq \emptyset$ of I. Then we have for any nontrivial $\vartheta \in \mathcal{D}(\mathcal{J})$

$$|\vartheta e(s\mathbf{v}\cdot\mathbf{\psi})|_{\mathbf{A}} = |\vartheta|_{\mathbf{A}} \neq 0$$
 for all $s \in \mathbf{R}$,

since $e(sv \cdot \psi)$ is the product of a unimodular complex number and a unitary character of **R**.

Thus (i) is not fulfilled, q.e.d.

Remark. — Condition (ii) of Corollary 1 is clearly satisfied if $\tau^{-1}(\{0\})$ has empty interior. As will be shown later (Lemma 3), this is always the case if ψ_1, \ldots, ψ_n are real analytic and linearly independent modulo affine mappings. However one should notice that global linear independence does not in general imply local linear independence.

Theorem 2.
$$-$$
 (i) If $\tau^{-1}(\{0\}) = \emptyset$, then for $\vartheta \in \mathscr{D}(I)$
$$|\vartheta e(\eta \cdot \psi)|_{PM} = 0(|\eta|^{-1/(n+1)}) \quad as \quad |\eta| \to \infty.$$

(ii) If $\vartheta \in \mathscr{D}(I)$, and if there exists an $x_0 \in I$ with $\vartheta(x_0) \neq 0$ and $\tau(x_0) \neq 0$, then there exists an $\varepsilon > 0$ and a function $\xi \in C^{\infty}((-\varepsilon, \varepsilon), \mathbf{R}^n)$ with

$$\det (\xi(y)\xi'(y) \dots \xi^{(n-1)}(y)) \neq 0 \quad \text{for all} \quad y \in (-\varepsilon, \varepsilon),$$

such that, for some C > 0,

$$|\Re(s\xi(v)\cdot\psi)|_{\rm DM} \ge C(1+|s|)^{-1/(n+1)}$$

for all $s \in \mathbf{R}$ and $y \in (-\varepsilon, \varepsilon)$.

Assume that $\tau^{-1}(\{0\})$ has empty interior. Then we have

COROLLARY 2. — There exists a $\vartheta \in \mathscr{D}(I)$, $\vartheta \neq 0$, such that for all positive $\alpha_1, \ldots, \alpha_n \in \mathbf{R}$ with $\sum_{i=1}^n \alpha_i \leqslant (n+1)^{-1}$, there exists a constant $C = C(\alpha_1, \ldots, \alpha_n) > 0$ such that

(2.1)
$$|\vartheta e(\eta \cdot \psi)|_{PM} \leqslant C \prod_{j=1}^{n} |\eta_{j}|^{-\alpha_{j}}.$$

Conversely, if $\alpha_1, \ldots, \alpha_n \in \mathbb{R}$ are positive, and if there exists $a \ \vartheta \in \mathcal{D}(I)$, $\vartheta \neq 0$, and $a \ C > 0$ such that (2.1) holds, then

$$\sum_{1}^{n}\alpha_{j}\leqslant (n+1)^{-1}.$$

Proof of Corollary 2. — If $\tau^{-1}(\{0\})$ has empty interior, then there is of course an $x_0 \in I$ with $\tau(x_0) \neq 0$, and so, for $\vartheta \in \mathcal{D}(I)$ with sufficiently small support near x_0 ,

$$|\vartheta e(\eta \cdot \psi)|_{PM} \leqslant C(1+|\eta|)^{-1/(n+1)}$$

by Theorem 2, (i).

If $\alpha_1, \ldots, \alpha_n$ are positive and $\Sigma \alpha_i \leq (n+1)^{-1}$, then

$$\prod_{j} |\eta_{j}|^{\alpha_{j}} \leqslant |\eta|^{1/(n+1)} \qquad \text{for} \qquad |\eta| \geqslant 1 \,,$$

hence

$$|\vartheta e(\eta \cdot \psi)|_{PM} \leqslant C \prod_{j} |\eta_{j}|^{-\alpha_{j}} \quad \text{for} \quad |\eta| \geqslant 1,$$

and the same estimate holds for all η if one replaces C by C + $|\vartheta|_{\Gamma^{\perp}}$.

Conversely, let now $\vartheta \in \mathcal{D}(I)$, $\vartheta \neq 0$, such that (2.1) holds for some $\alpha_i \geq 0$, and assume

$$\Sigma \alpha_i = (n+1)^{-1} + \delta, \qquad \delta > 0.$$

Since $\tau^{-1}(\{0\})$ has empty interior, there is an $x_0 \in I$ with $\vartheta(x_0) \neq 0$ and $\tau(x_0) \neq 0$. Choose $\varepsilon > 0$ and $\xi \in C^{\infty}((-\varepsilon, \varepsilon), \mathbb{R}^n)$ as in Theorem 2 (ii). Since $\det(\xi(y)\xi'(y)\ldots\xi^{(n-1)}(y)) \neq 0$ for all $y \in (-\varepsilon, \varepsilon)$, there exists a $y_0 \in (-\varepsilon, \varepsilon)$ with

$$\xi_j(y_0) \neq 0$$
 for $j = 1, \ldots, n$.

It follows

$$|\vartheta e(s\xi(y_0)\cdot\psi)|_{PM}\geqslant C'(1+|s|)^{-1/(n+1)}.$$

On the other hand, (2.1) yields

$$\begin{split} |\Im e(s\xi(y_0)\cdot \psi)|_{\mathsf{PM}} &\leqslant C \prod_j |s\xi_j(y_0)|^{-\alpha_j} \\ &= \left(C \prod_j |\xi_j(y_0)|^{-\alpha_j}\right) |s|^{-1/(n+1)} |s|^{-\delta} \,. \end{split}$$

For |s| sufficiently large this leads to a contradiction to (2.2), q.e.d.

Corollary 2 demonstrates that the result in Theorem 2 is in some sense best possible.

3.

Before we start to prove the theorems above we will state some lemmas. The first one is due to J.-E. Björk and is cited in [3], Lemma 1.6:

LEMMA 1. – Let $I \neq \emptyset$ be a bounded, open interval in \mathbb{R} , and let $\varphi \in \mathcal{D}(I)$, $g \in \mathbb{C}^p(I)$ with

$$0 < C_1 \le |g'(x)| + |g''(x)| + \cdots + |g^{(p)}(x)| \le C_2$$

if $x \in \overline{I}$, where C_1 and C_2 are constants and p is a positive integer. Then there exists a constant C not depending on g, such that

$$\left| \int \varphi(x) e^{2\pi i t g(x)} dx \right| \leqslant C(1+|t|)^{-1/p}$$

for every $t \in \mathbf{R}$.

The second lemma will be used to prove the remark following Corollary 1. I would like to thank Professor H. Leptin for pointing out to me a shorter proof than my original one. By $\langle A \rangle$ we denote the exterior product in the Grassmann algebra $\Lambda(\mathbb{R}^n)$.

LEMMA 2. – Let $\psi \in C^{\infty}(I, \mathbb{R}^n)$. Then

$$\psi(x) \wedge \psi'(x) \ldots \wedge \psi^{(n-1)}(x) = 0$$

for all $x \in I$ implies

$$\psi^{(k_1)}(x) \wedge \psi^{(k_2)}(x) \wedge \ldots \wedge \psi^{(k_n)}(x) = 0$$

for all $x \in I$ and $k_1, \ldots, k_n \in \mathbb{N}_0$.

Proof. – Fix $x_0 \in I$, and assume first $\psi(x_0) \neq 0$. If $u \in C^{\infty}(I, \mathbb{R})$, then

$$(u\psi)^{(k)} = \sum_{j=0}^{k} {k \choose j} u^{(k-j)} \psi^{(j)},$$

so $\psi \wedge \psi' \wedge \ldots \wedge \psi^{(n-1)} \equiv 0$ implies

$$(u\psi) \wedge (u\psi)' \wedge \ldots \wedge (u\psi)^{(n-1)} \equiv 0.$$

So, it is no loss of generality to assume

$$\psi_n(x) = 1$$
 for $x \in I$.

If $\{e_j\}_j$ denotes the canonical basis of \mathbf{R}^n , we may thus write $\psi(x) = \sum_{j=1}^{n-1} \psi_j(x) e_j + e_n = \rho(x) + e_n$, where $\rho(x) \in \mathbf{R}^{n-1} \times \{0\} \subset \mathbf{R}^n$. This yields

$$0 = \psi(x) \wedge \psi'(x) \wedge \ldots \wedge \psi^{n-1}(x) = \rho(x) \wedge \rho'(x) \wedge \ldots \wedge \rho^{(n-1)}(x) + e_n \wedge \rho'(x) \wedge \ldots \wedge \rho^{(n-1)}(x),$$

and since $\rho(x)$, $\rho'(x)$, ..., $\rho^{(n-1)}(x)$ are clearly linearly dependent, we get

$$0 = \rho'(x) \wedge \rho''(x) \wedge \ldots \wedge \rho^{(n-1)}(x).$$

By induction over n, we now may assume

$$0 = \rho^{(k_2)}(x) \wedge \rho^{(k_3)}(x) \wedge \ldots \wedge \rho^{(k_n)}(x)$$

for $x \in I$ and $k_i \ge 1$.

This implies

$$\psi^{(k_1)}(x) \wedge \ldots \wedge \psi^{(k_n)}(x) = e_n^{(k_1)}(x) \wedge \rho^{(k_2)}(x) \wedge \ldots \wedge \rho^{(k_n)}(x) = 0$$

for $0 \le k_1 < k_2 < \cdots < k_n$, where we considered e_n as the function $e_n(x) = e_n$.

Thus we have proved

$$\psi^{(k_1)}(x_0) \wedge \psi^{(k_2)}(x_0) \wedge \ldots \wedge \psi^{(k_n)}(x_0) = 0$$

for all $x_0 \in I_0 = \{x \in I : \psi(x) \neq 0\}$ and $k_j \geqslant 0$. By continuity, the same holds true for $x_0 \in \overline{I}_0 \land I$, hence for all $x_0 \in I$, since for $y \in I \setminus \overline{I}_0$ clearly $\psi^{(k)}(y) = 0$ for every $k \in \mathbb{N}_0$.

LEMMA 3. – If $\psi = (\psi_1, \ldots, \psi_n) \in C^{\infty}(I, \mathbb{R}^n)$ is real analytic, and if ψ_1, \ldots, ψ_n are linearly independent modulo affine mappings, then $\tau^{-1}(\{0\})$ has empty interior, where τ denotes the torsion of ψ .

Proof. Assume $\tau(x) = 0$ for every x in some nonempty open interval $J \subset I$. Fix $x_0 \in J$. Then, passing to a possibly smaller interval, we may assume that ψ_i has an absolute convergent series expansion

$$\psi_j(x) = \sum_{k=0}^{\infty} a_k^j (x - x_0)^k, \quad j = 1, \ldots, n, \quad x \in J.$$

Define vectors

$$a_k = (a_k^j)_{j=1,\ldots,n} \in \mathbf{R}^n$$

and

$$a^j = (a_k^j)_{k=2,\ldots,\infty} \in \mathbb{R}^{N_1}, \qquad \mathbb{N}_1 = \mathbb{N} \setminus \{0,1\}.$$

By Lemma 2, $\psi^{(k_1)}(x_0), \ldots, \psi^{(k_n)}(x_0)$ are linearly dependent for any $k_j \in \mathbb{N}$ with $2 \leq k_1 < \ldots < k_n$, i.e. a_{k_1}, \ldots, a_{k_n} are linearly dependent for $2 \leq k_1 < \ldots < k_n$. But this implies that a^1, \ldots, a^n are linearly

dependent, i.e. there exist $v_1, \ldots, v_n \in \mathbb{R}$, not all zero, with

$$0 = \sum_j v_j a^j, \quad \text{i.e.}$$

$$\sum_j v_j \psi_j(x) = \sum_j v_j a^j_0 + v_j a^j_1(x - x_0) \qquad \text{for} \qquad x \in J.$$

But, since ψ is real analytic, this equation holds for all $x \in I$, i.e. $\sum_{j} v_{j} \psi_{j}$ is affine linear.

4.

Proof of Theorem 1. — It is well-known (see e.g. [1], [7]) that for $\varphi \in \mathcal{D}(\mathbf{R})$ one has the estimate

$$|\varphi|_{A} \leqslant \left\{2 \left| \sup \varphi \right| |\varphi|_{\infty} |\varphi'|_{\infty}\right\}^{1/2},$$

where $|\text{supp }\phi|$ denotes the Lebesgue measure of the support of ϕ . From (4.1) one immediately gets (i) of Theorem 1.

Now, suppose there exists a subinterval J in I and a $\sigma > 0$ such that $|\vartheta(x)| \ge \sigma$ and $|\vartheta(x) - \vartheta(y)| < \sigma/2$ for $x, y \in J$, and such that ψ_1, \ldots, ψ_n are linearly independent modulo affine mappings on J. Then a simple compactness argument yields:

There are constants $\varepsilon > 0$, $\delta > 0$, such that for every $\eta \in \mathbb{R}^n$ with $|\eta| = 1$ there is an interval J_{η} of length 2ε in J with

$$(4.2) |\eta \cdot \psi''(x)| \ge \delta \text{for all} x \in J_{\eta}.$$

Now choose $\varphi \in \mathcal{D}(-\varepsilon,\varepsilon)$, $\varphi \geqslant 0$, with $\int \varphi(x) dx = 1$. For fixed $\eta \in \mathbb{R}^n$, $\eta \neq 0$, set $\eta' = |\eta|^{-1} \eta$, and choose $J_{\eta'}$ as in (4.2). Let $\tilde{\varphi}$ be a suitable translate of φ such that supp $\tilde{\varphi} \subset J_{\eta'}$. Then we get

$$(4.3) \quad 0 < \sigma/2 \le \left| \int \vartheta(x) \tilde{\varphi}(x) \, dx \right|$$

$$= \left| \int \vartheta(x) e(\eta \cdot \psi)(x) \tilde{\varphi}(x) e(-\eta \cdot \psi)(x) \, dx \right|$$

$$\le \left| \vartheta e(\eta \cdot \psi) \right|_{A} \left| \tilde{\varphi} e(-\eta \cdot \psi) \right|_{PM},$$

since $J_{\eta'} \subset J$.

For $\xi \in \mathbf{R}$ one has

$$\{\tilde{\varphi}e(\eta \cdot \psi)\}^{\hat{}}(-\xi) = \int \tilde{\varphi}(x)e(-\xi x - \eta \cdot \psi(x)) dx$$
$$= \int \varphi(x)e(-|\eta|g(x)) dx,$$

where g is a function on $[-\varepsilon,\varepsilon]$ which is a certain translate of the function

$$x \mapsto \xi' x + \eta' \cdot \psi(x)$$
 on $J_{\eta'}$,

where $\xi' = |\eta|^{-1} \xi$.

But (4.2) implies

$$\delta \leq |g''(x)|$$
 for every $x \in [-\varepsilon, \varepsilon]$.

Moreover, if we set $A = 2 \sup_{x \in J} |\psi'(x)|$, $B = \sup_{x \in J} |\psi''(x)|$, then for $|\xi| \le A|\eta|$:

$$|g'(x)| + |g''(x)| \le |\xi'| + |\eta'|(A+B)$$

 $\le 2A + B$

for every $x \in [-\varepsilon, \varepsilon]$.

Thus, by Lemma 1, there exists a C>0, such that for $|\xi|\leqslant A|\eta|$

$$(4.4) \qquad \left| \int \tilde{\varphi}(x) e(-\xi x - \eta \cdot \psi(x)) \ dx \right| \leq C(1 + |\eta|)^{-1/2}.$$

And, if $|\xi| > A|\eta|$, then integration by parts yields

(4.5)
$$\left| \int \tilde{\varphi}(x) e(-\xi x - \eta \cdot \psi(x)) \, dx \right| = \left| \int e(-|\eta|g(x)) \left(\frac{\varphi}{2\pi i |\eta|g'} \right)'(x) \, dx \right| \\ \leq (2\pi |\eta|)^{-1} \int \left\{ \frac{|\varphi'(x)|}{|g'(x)|} + \frac{|\varphi(x)| |g''(x)|}{|g'(x)|^2} \right\} dx \\ \leq C' |\eta|^{-1},$$

where C' is some constant depending on φ , ψ and A only, since for $x \in [-\varepsilon, \varepsilon]$ we have $|g''(x)| \leq B$ and $|g'(x)| = |\xi' + \eta' \psi'(y)| \geq A - A/2$ for some $y \in J$.

Now, by (4.4), (4.5),

$$|\tilde{\varphi}e(-\eta\cdot\psi)|_{PM} \leqslant (C+C')|\eta|^{-1/2}$$
 if $|\eta| \geqslant 1$,

which together with (4.3) proves Theorem 1 (ii).

Proof of Theorem 2. — Assume $\tau(x) \neq 0$ for every $x \in I$, and let $\theta \in \mathcal{D}(I)$, $\theta \neq 0$. Passing to a smaller interval, we may even assume that I is closed.

Set $A = 2 \sup_{x \in I} |\psi'(x)|$, and for $\xi' \in \mathbb{R}$, $|\xi'| \le A$, $\eta' \in \mathbb{R}^n$, $|\eta'| = 1$, $x \in I$ let

$$Q_{\xi',\eta'}(x) = \sum_{j=1}^{n+1} |(\xi'x + \eta' \cdot \psi(x))^{(j)}(x)|.$$

Since $\tau^{-1}(\{0\}) = \emptyset$, we have $Q_{\xi',\eta'}(x) \neq 0$ for every $x \in I$, and since $Q_{\xi',\eta'}(x)$ is continuous in ξ',η' and x on the compact space $[-A,A] \times \{\eta' \in \mathbf{R}^n : |\eta'| = 1\} \times I$, there exist constants $C_1 > 0$, $C_2 > 0$, such that

$$(4.6) C_1 \leqslant Q_{\xi',\eta'}(x) \leqslant C_2$$

for all $x \in I$, ξ' , η' with $|\xi'| \le A$, $|\eta'| = 1$.

So, using quite the same arguments as in the proof of Theorem 1 (ii), we can deduce from (4.6) by Lemma 1:

$$|\vartheta e(\eta \cdot \psi)|_{PM} \leqslant C(1+|\eta|)^{-1/(n+1)}$$

for some constant C > 0, which proves (i).

To prove (ii), we will assume, for convenience, $x_0 = 0$, i.e. $0 \in I$, and $\vartheta(0) \neq 0$, $\tau(0) \neq 0$.

Let $\varepsilon > 0$ such that $\tau(x) \neq 0$ for $x \in [-\varepsilon, \varepsilon]$.

Since $\psi''(x)$, $\psi'''(x)$, ..., $\psi^{(n+1)}(x)$ are linearly independent for $x \in [-\varepsilon, \varepsilon]$, there exists a function $\xi \in C^{\infty}([-\varepsilon, \varepsilon], \mathbb{R}^n)$, such that for every $x \in [-\varepsilon, \varepsilon]$

(4.7)
$$\xi(x) \cdot \psi^{(j)}(x) = 0, \quad j = 2, \ldots, n,$$

and

(4.8)
$$\xi(x) \cdot \hat{\Psi}^{(n+1)}(x) = 1.$$

Differentiating (4.7) and inserting (4.8), we get

$$\xi'(x) \cdot \psi^{(j)}(x) = 0$$
 for $j = 2, ..., n-1$,

and

$$\xi'(x)\psi^{(n)}(x) = -1.$$

Repeating this process, one inductively obtains for k = 0, ..., n - 1

(4.9)
$$\begin{cases} \xi^{(k)}(x) \cdot \psi^{(j)}(x) = 0 & \text{for } j = 2, \dots, n-k, \\ \xi^{(k)}(x) \cdot \psi^{(n+1-k)}(x) = (-1)^k. \end{cases}$$

So, if we define matrices

$$S(x) = (\xi_j^{(n-i)}(x))_{i,j=1,...,n}, \qquad T(x) = (\psi_i^{(j+1)}(x))_{i,j=1,...,n},$$

then (4.9) means that S(x)T(x) is an upper triangular matrix with diagonal elements 1 or -1, which yields

$$(4.10) |\det (\xi(x)\xi'(x) \dots \xi^{(n-1)}(x))| = |\det S(x)| = |\tau(x)|^{-1} \neq 0$$

for all $x \in [-\varepsilon, \varepsilon]$.

We now claim:

There is a constant C > 0, such that for all $y \in (-\varepsilon, \varepsilon)$ and $s \in \mathbb{R}$

(4.11)
$$|\vartheta e(s\xi(y)\cdot\psi)|_{PM} \geqslant C(1+|s|)^{-1/(n+1)}.$$

Choose $y \in (-\varepsilon, \varepsilon)$. Then by (4.7), $(\xi(y) \cdot \psi)^{(j)}(y) = \delta_{j,n+1}$ for j = 2, ..., n+1, and so a Taylor expansion of $\xi(y) \cdot \psi$ yields (for ε small enough)

(4.12)
$$(\xi(y) \cdot \psi)(x) = \alpha + \beta x + (x - y)^{n+1} g(x)$$
 for $x \in (-2\varepsilon, 2\varepsilon)$,

where g is some smooth function on $(-2\varepsilon,2\varepsilon)$ which depends on y, and where α and β are some real numbers.

Let us remark here that although $g = g_y$ depends on y, $\sup_{|x|<2\varepsilon} |g_y'(x)|$ is uniformly bounded for $y \in (-\varepsilon, \varepsilon)$.

Now take $\rho \in \mathcal{D}(\mathbf{R})$ with supp $\rho \subset (-\varepsilon, \varepsilon)$, $\rho \geqslant 0$ and $\int \rho(x) dx = 1$, and set $\tilde{\rho}(x) = \rho(|s|^{1/(n+1)}(x-y))$.

If we choose ε small enough such that

$$|\vartheta(0)-\vartheta(x)|<\frac{1}{2}\,|\vartheta(0)|$$

for $x \in (-2\varepsilon, 2\varepsilon)$, then we get

$$\left| \int \vartheta(x) \tilde{\rho}(x) \, dx \right| = \left| \int \vartheta(|s|^{-1/(n+1)}x + y) \rho(x) \, dx \right| |s|^{-1/(n+1)}$$

$$\geq \frac{1}{2} |\vartheta(0)| |s|^{-1/(n+1)}, \quad \text{if} \quad |s| \geq 1;$$

and since

$$\left| \int \vartheta(x) \, \tilde{\rho}(x) \, dx \right| = \left| \int \vartheta(x) e(s\xi(y) \cdot \psi) \tilde{\rho}(x) e(-s\xi(y) \cdot \psi) \, dx \right|$$

$$\leq \left| \vartheta e(s\xi(y) \cdot \psi)_{PM} |\tilde{\rho}e(-s\xi(y) \cdot \psi)|_{A},$$

(4.11) will follow if we can show that $|\tilde{\rho}e(-s\xi(y)\cdot\psi)|_A$ is uniformly bounded for $y \in (-\varepsilon,\varepsilon)$ and $|s| \ge 1$.

Now, regular affine mappings of R induce isometries of the Fourier algebra A = A(R), thus

$$|\tilde{\rho}e(-s\xi(y)\cdot\psi)|_{A} = |\rho e(-s\xi(y)\cdot\widetilde{\psi}|_{A},$$

where $\tilde{\psi}(x) = \psi(|s|^{-1/(n+1)}x + y)$.

Since for $x \in \text{supp } \rho$ and $|s| \ge 1$,

$$|s|^{-1/(n+1)}x + y \in (-2\varepsilon, 2\varepsilon),$$

(4.12) yields

$$\xi(y) \cdot \widetilde{\psi}(x) = \alpha + \beta y + \beta |s|^{-1/(n+1)} x + |s|^{-1} x^{n+1} g(|s|^{-1/(n+1)} x + y).$$

Thus

$$|\tilde{\rho}e(-s\xi(y)\cdot\psi)|_{A}=|\rho e(h)|_{A}$$

where $h(x) = -s|s|^{-1}x^{n+1}g(|s|^{-1/(n+1)}x+y)$. If we again apply estimate (4.1), we easily see that $|pe(h)|_A$ is uniformly bounded for $y \in (-\varepsilon,\varepsilon)$ and $|s| \ge 1$, q.e.d.

BIBLIOGRAPHY

- [1] F. CARLSON, Une inégalité, Ark. Mat. Astr. Fys., 25, B1 (1934).
- [2] L. Corwin, F. P. Greenleaf, Singular Fourier integral operators and representations of nilpotent Lie groups, *Comm. on Pure and Applied Math.*, B1 (1978), 681-705.
- [3] Y. DOMAR, On the Banach algebra $A(\Gamma)$ for smooth sets $\Gamma \subset \mathbb{R}^n$, Comment. Math. Helv., 52 (1977), 357-371.
- [4] C. S. Herz, Fourier transforms related to convex sets, *Ann. of Math.*, (2), 75 (1962), 215-254.
- [5] L. HÖRMANDER, Lower bounds at infinity for solutions of differential equations with constant coefficients, *Israel J. Math.*, 16 (1973), 103-116.
- [6] W. LITTMAN, Fourier transforms of surface-carried measures and differentiability of surface averages, Bull, Amer. Math. Soc., 69 (1963), 766-770.
- [7] D. Müller, On the spectral synthesis problem for hypersurfaces of \mathbb{R}^n , J. Functional Analysis, 47 (1982), 247-280.

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