# On the explicit calculation of fundamental solutions 

Peter Wagner<br>University of Innsbruck, Technikerstr. 13, A-6020 Innsbruck, Austria<br>Received 24 October 2003<br>Submitted by R.M. Aron<br>Dedicated to Professor J. Horváth on the occasion of his 80th birthday

## Introduction

The notion of a fundamental solution (in the sequel abbreviated as FS) gradually became clearer during the 19th and 20th century. It is only in the setting of Laurent Schwartz' theory of distributions that FSs can be defined in general and can be applied-via the convolution of distributions-to the solution of linear partial differential equations with constant coefficients. In part, the relevant concepts were worked out by John Horváth, cf. [20-25]. In this survey paper, we first review some important steps in the "history of FSs." Second, we explain why the singular support of the FSs of homogeneous operators $P(\partial)$ is just the dual hypersurface of the zero set of $P$ if $P$ is of principal type. This is of fundamental importance in the third part, where we present some recent results in the calculation of FSs of homogeneous cubic and quartic operators in three dimensions. Finally, we discuss what is known for the system of crystal optics, where, similarly as in dynamic anisotropic elasticity, many questions are still open.

## 1. A brief history of fundamental solutions

### 1.1. Fundamental solutions in the 18 th and 19th century: special equations of mathematical physics

The first use of a non-trivial fundamental solution can probably be ascribed to Jean d'Alembert. In 1747 he considered the deflection $u$ of a vibrating string. It satisfies the

[^0]one-dimensional wave equation
$$
\frac{\partial^{2} u}{\partial t^{2}}-c^{2} \frac{\partial^{2} u}{\partial x^{2}}=f
$$
which is solved by convolving $f$ with the $\operatorname{FS} E(t, x)=\frac{1}{2 c} Y(t-|x| / c)$ of the operator $\partial_{t}^{2}-c^{2} \partial_{x}^{2}, c>0$. In fact, this yields the formula
$$
u(t, x)=f * E=\frac{1}{2 c} \iint_{|x-\xi|<c(t-\tau)} f(\tau, \xi) \mathrm{d} \tau \mathrm{~d} \xi
$$
which—applied to the initial value problem, i.e., with $f=\delta(t) \otimes u_{1}(x)+\delta^{\prime}(t) \otimes$ $u_{0}(x)$, where $u_{j}(x)=\left(\partial_{t}^{j} u\right)(0, x), j=0,1$,—furnishes d'Alembert's solution, see [1], [33, p. 15 ff$]$. We observe that the FS $E$ appears, as in much of the old literature, only in an implicit way.

Here and in the following, we use the notations $\partial_{t}=\frac{\partial}{\partial t}, \partial_{x}=\frac{\partial}{\partial x}, \partial_{1}=\frac{\partial}{\partial x_{1}}$ etc., $\Delta_{n}=$ $\partial_{1}^{2}+\cdots+\partial_{n}^{2}, \partial^{\alpha}=\partial_{1}^{\alpha_{1}} \cdots \partial_{n}^{\alpha_{n}}, P(\partial)=\sum_{\alpha \in \mathbb{N}_{0}^{n}} c_{\alpha} \partial^{\alpha}, Y(t)=1$ for $t \geqslant 0$ and $Y(t)=0$ else, $x=\left(x_{1}, \ldots, x_{n}\right)^{\mathrm{T}}, \nabla=\left(\partial_{1}, \ldots, \partial_{n}\right)^{\mathrm{T}},|x|=\sqrt{x_{1}^{2}+\cdots+x_{n}^{2}}$.

In 1789, Pierre Simon de Laplace used the FS $E=-\frac{1}{4 \pi|x|}$ of the elliptic operator $\Delta_{3}$, which bears his name, and thereby established the connexion of the Laplace operator with the Newtonian gravitational potential (cf. [29]). To tell the truth, Laplace just recognized that $\Delta_{3}(E * f)=0$ outside the support of $f$, and it was Simon Denis Poisson, who obtained the equation $\Delta_{3}(E * f)=f$ in 1813 (cf. [41]).

In 1809, Laplace considered the first parabolic operator, namely the heat operator $\partial_{t}-\Delta_{n}$, and calculated its FS

$$
E(t, x)=\frac{Y(t)}{(4 \pi t)^{n / 2}} \mathrm{e}^{-|x|^{2} /(4 t)}
$$

in the case $n=1$, cf. [30]. The generalization to higher $n$, in particular to $n=2$, was found by Poisson in 1818 [42].

In 1818 , Joseph Fourier was able to calculate the FS $E$ of the operator of the dynamic deflections of beams $\partial_{t}^{2}+\partial_{x}^{4}$, an operator of fourth order:

$$
\begin{aligned}
E(t, x) & =\frac{Y(t)}{2 \sqrt{\pi}} \int_{0}^{t} \sin \left(\frac{x^{2}}{4 \tau}+\frac{\pi}{4}\right) \frac{\mathrm{d} \tau}{\sqrt{\tau}} \\
& =Y(t)\left[\sqrt{\frac{t}{\pi}} \sin \left(\frac{x^{2}}{4 t}+\frac{\pi}{4}\right)-\frac{|x|}{2} C\left(\frac{x^{2}}{4 t}\right)+\frac{|x|}{2} S\left(\frac{x^{2}}{4 t}\right)\right],
\end{aligned}
$$

where

$$
\left.\begin{array}{c}
C(x) \\
S(x)
\end{array}\right\}=\frac{1}{\sqrt{2 \pi}} \int_{0}^{x}\left\{\begin{array}{c}
\cos \\
\sin
\end{array}\right\}(u) \frac{\mathrm{d} u}{\sqrt{u}},
$$

see [7].

As well in 1818, Poisson generalized d'Alembert's formula to three space dimensions by representing the solutions of the wave operator $\partial_{t}^{2}-\Delta_{3}$ as convolution with the FS $E=\delta(t-|x|) /(4 \pi|x|)$, cf. [43]. This notation, viz. the first use of Dirac's delta function, goes back to Gustav Kirchhoff's paper of 1882 (see [27], [33, p. 99]).

In 1849, George Stokes obtained-as the kernel of an integral representation-the fundamental matrix $E$ of the system of partial differential operators which describes elastic waves in isotropic media [46]. This system can be found already in a memoir of 1829 by Poisson (cf. [44]). It is given by

$$
P(\partial)=\left(\partial_{t}^{2}-\mu \Delta_{3}\right) I_{3}-(\lambda+\mu) \nabla \cdot \nabla^{\mathrm{T}} \quad\left(I_{3}=3 \text { by } 3 \text { unit matrix }\right)
$$

and Stokes' fundamental matrix reads

$$
\begin{aligned}
E(t, x)= & \frac{I_{3}|x|^{2}-x \cdot x^{\mathrm{T}}}{4 \pi \mu|x|^{3}} \delta\left(t-\frac{|x|}{\sqrt{\mu}}\right)+\frac{x \cdot x^{\mathrm{T}}}{4 \pi(\lambda+2 \mu)|x|^{3}} \delta\left(t-\frac{|x|}{\sqrt{\lambda+2 \mu}}\right) \\
& +\frac{t}{4 \pi|x|^{3}}\left(I_{3}-\frac{3 x \cdot x^{\mathrm{T}}}{|x|^{2}}\right)\left[Y\left(t-\frac{|x|}{\sqrt{\mu}}\right)-Y\left(t-\frac{|x|}{\sqrt{\lambda+2 \mu}}\right)\right] .
\end{aligned}
$$

The FS $E=Y(t-|x|) /\left(2 \pi \sqrt{t^{2}-|x|^{2}}\right)$ of the wave operator in two space dimensions, i.e., of $\partial_{t}^{2}-\Delta_{2}$, was found as late as 1894 by Vito Volterra, cf. [52].

### 1.2. Fundamental solutions in the 20th century: general theories

Investigating the equations of static anisotropic elasticity, Ivar Fredholm found in 1900 (cf. [8]) the fundamental matrix $E$ of the elliptic 3 by 3 system

$$
P(\partial)=\left(\sum_{k, l=1}^{3} c_{i j k l} \partial_{k} \partial_{l}\right)_{i, j=1,2,3}, \quad c_{i j k l} \in \mathbf{R},
$$

of linear partial differential operators in three variables with constant coefficients and homogeneous of second order. In our notation, his result is the following (cf. [8, (10), p. 7], [39, 3.2.2, (F), p. 332]):

$$
E(x)=-\frac{i \operatorname{sign}\left(x_{2}\right)}{2 \pi} \sum_{k=1}^{3} \frac{\left|\zeta_{k}(x)\right|^{2} P\left(\zeta_{k}(x)\right)^{\mathrm{ad}}}{x_{2} \frac{\partial \operatorname{det} P}{\partial \xi_{1}}\left(\zeta_{k}(x)\right)-x_{1} \frac{\partial \operatorname{det} P}{\partial \xi_{2}}\left(\zeta_{k}(x)\right)},
$$

where $P(\zeta)^{\text {ad }}$ denotes the adjoint matrix of $P(\zeta)$ and $\zeta_{k}(x) \in \mathbf{C}^{3} \backslash\{0\}$ are determined up to complex factors by the conditions

$$
x \cdot \zeta_{k}(x)=0, \quad \operatorname{det} P\left(\zeta_{k}(x)\right)=0, \quad \operatorname{Im}\left(\frac{\zeta_{k}(x)_{1}}{\zeta_{k}(x)_{3}}\right)>0, \quad k=1,2,3 .
$$

In 1908, Fredholm succeeded in representing the FSs of elliptic homogeneous operators in 3 variables by Abelian integrals [9]. To test the theory, he applied it to the operator $\partial_{1}^{4}+\partial_{2}^{4}+\partial_{3}^{4}$, and he obtained, up to the constant factor $-\frac{1}{8 \pi}$, the beautiful formula

$$
E(x)=-\frac{1}{8 \pi} \sum_{j=1}^{3}\left|x_{j}\right| \int_{\zeta /\left(2 x_{j}^{2}\right)}^{\infty} \frac{\mathrm{d} u}{\sqrt{4 u^{3}-u}}
$$

$$
=-\frac{1}{8 \pi} \sum_{j=1}^{3} x_{j} F\left(\arcsin \left(\frac{\sqrt{2} x_{j}}{\sqrt{\zeta+x_{j}^{2}}}\right), \frac{1}{\sqrt{2}}\right)
$$

Therein $\zeta$ denotes the largest of the three real roots of the cubic

$$
\zeta^{3}-\left(x_{1}^{4}+x_{2}^{4}+x_{3}^{4}\right) \zeta-2 x_{1}^{2} x_{2}^{2} x_{3}^{2}=0
$$

and $F$ denotes the elliptic integral of the first kind, cf. [12, 3.131.8 and 8.111]. (A generalization to elliptic operators of the form $\sum_{j, k=1}^{3} c_{j k} \partial_{j}^{2} \partial_{k}^{2}$ can be found in [56].)

In 1911, his pupil Nils Zeilon gave the first definition of a FS in case it is a locally integrable function (cf. [58]):
$F$ is a FS of $f\left(\frac{\partial}{\partial x}, \frac{\partial}{\partial y}, \frac{\partial}{\partial z}\right)$ if and only if

$$
u=\iiint F(x-\lambda, y-\mu, z-v) \phi(\lambda, \mu, v) \mathrm{d} \lambda \mathrm{~d} \mu \mathrm{~d} v
$$

solves $f\left(\frac{\partial}{\partial x}, \frac{\partial}{\partial y}, \frac{\partial}{\partial z}\right) u=\phi$.
In 1913, Zeilon transferred Fredholm's results to non-elliptic operators and, in particular, he considered $\partial_{1}^{3}+\partial_{2}^{3}+\partial_{3}^{3}$ (cf. [59]). He determined the singular support of its FS, but, in contrast to Fredholm, he was not able to obtain an explicit representation for the FS. However, explicit formulae were found recently, see [53-55].

In three famous papers from 1926 to 1928 (see [13]), Gustav Herglotz overcame the restriction to 2 or 3 independent variables and represented the FSs of elliptic and of strictly hyperbolic homogeneous operators of the degree $m$ in $n$ variables (with $n \leqslant m$ ) by ( $n-1$ )fold and by $(n-2)$-fold integrals, respectively. Later, these formulae came to be known as the Herglotz-Petrovsky formulae.

In 1945, Ivan Petrovsky represented-in the hyperbolic case-the FS $E$ by integrals over cycles in complex projective space and investigated the lacunas of $E$ by means of algebraic topology [40].

In 1950/51, Laurent Schwartz first published his Théorie des Distributions [45], in which framework he also gave the general definition of FSs:

$$
E \text { is a FS of } P(\partial) \quad \text { if and only if } \quad P(\partial) E=\delta
$$

In Chapter 6 (Transformation de Fourier) of his book, Schwartz rederives the FSs of ( $\Delta_{n}-$ $\lambda)^{m},\left(\partial_{t}^{2}-\Delta_{n}-\lambda\right)^{m},\left(\partial_{t}-\Delta_{n}-\lambda\right)^{m}, \lambda \in \mathbf{C}$, by distributional calculus, cf. also [22,23].

In 1952, Jean Leray stated a distributional version of the Herglotz-Petrovsky formulae for homogeneous hyperbolic operators, thereby also treating the case $m<n$ (cf. [31]). The same goal was reached in 1959 by Vladimir A. Borovikov for operators of principal type (cf. [5]) and presented in the textbook "Generalized Functions" by Israel M. Gel'fand and Georgi E. Shilov (cf. [11]).

The first existence proofs for FSs were given in 1953/54 by Bernard Malgrange and Leon Ehrenpreis (cf. [6,34]). These proofs were based on the Hahn-Banach theorem. In 1957, Lars Hörmander showed that there always exist "regular" FSs (at that time called "proper" FSs) having "best" regularity properties (cf. [15]). The existence of FSs depending $\mathcal{C}^{\infty}$ or even holomorphic (in case of "constant strength") on the coefficients of $P(\partial)$ was proved by François Trèves, cf. [47-49]; see also the survey paper [38].

In 1957/58, Lars Hörmander and Stanislaw Łojasiewicz independently solved the "division problem" and thereby proved the existence of temperate FSs (cf. [16,32]). Different proofs thereof were found later by Michael F. Atiyah [2] and Joseph N. Bernstein [4].

In 1970/73, Michael Atiyah, Raoul Bott, and Lars Gårding extended and generalized Petrovsky's work, thereby developing a general theory of FSs of hyperbolic operators, cf. [3]. For general operators, this was established in the fundamental work of Lars Hörmander, cf. [14,17,18].

We also mention the first major table of FSs by Norbert Ortner in 1980 (cf. [36]) and the discovery of the connexion of lacunas of FSs with the existence of right inverses by Reinhold Meise, B. Alan Taylor, and Dietmar Vogt in 1990 (cf. [35]).

Finally, I would like to sketch a proof of the Malgrange-Ehrenpreis theorem I found in 1994, influenced by a paper of Heinz König (cf. [28]). This constructive proof seems to be the shortest one at present.

Theorem (Malgrange/Ehrenpreis, 1953/54). Let $P(\xi)=\sum_{|\alpha| \leqslant m} c_{\alpha} \xi^{\alpha}$ be a not identically vanishing polynomial in $\mathbf{R}^{n}$ (i.e., $c_{\alpha} \in \mathbf{C}, \xi=\left(\xi_{1}, \ldots, \xi_{n}\right) \in \mathbf{R}^{n}$, $\xi^{\alpha}=\xi_{1}^{\alpha_{1}} \cdots \xi_{n}^{\alpha_{n}}$, not all $\left.c_{\alpha}=0\right)$. Then there exists a FS of $P(\partial)$, i.e., $\exists E \in \mathcal{D}^{\prime}\left(\mathbf{R}^{n}\right): P(\partial) E=\delta$.

Proof [37]. The distribution $E \in \mathcal{D}^{\prime}\left(\mathbf{R}^{n}\right)$ defined by

$$
\begin{equation*}
E(x)=\frac{1}{\overline{P_{m}(\eta)}} \int_{\lambda \in \mathbb{C},|\lambda|=1} \lambda^{m} \mathrm{e}^{\lambda \eta x} \mathcal{F}^{-1}\left(\frac{\overline{P(\mathrm{i} \xi+\lambda \eta)}}{P(\mathrm{i} \xi+\lambda \eta)}\right) \frac{\mathrm{d} \lambda}{2 \pi \mathrm{i} \lambda} \tag{1}
\end{equation*}
$$

is a FS of $P(\partial)$, if $P_{m}(\xi)=\sum_{|\alpha|=m} c_{\alpha} \xi^{\alpha}$ (i.e., $P_{m}$ is the principal part of $P$ ), $\eta \in \mathbf{C}^{n}$ with $P_{m}(\eta) \neq 0$ is fixed, $\eta x=\eta_{1} x_{1}+\cdots+\eta_{n} x_{n}$, and $\mathcal{F}$ denotes the Fourier transform $\left((\mathcal{F} \phi)(x)=\int \phi(\xi) \mathrm{e}^{-\mathrm{i} x \xi} \mathrm{~d} \xi\right.$ for $\phi \in \mathcal{D}$, and extended to $\mathcal{S}^{\prime}$ by continuity) with the inverse $\mathcal{F}^{-1} T=(2 \pi)^{-n}(\mathcal{F} T)(-x)$. Formula (1) makes sense, since

$$
\frac{\overline{P(\mathrm{i} \xi+\lambda \eta)}}{P(\mathrm{i} \xi+\lambda \eta)} \in L^{\infty}\left(\mathbb{R}_{\xi}^{n}\right) \subset \mathcal{S}^{\prime}\left(\mathbb{R}^{n}\right)
$$

and since this distribution continuously depends on $\lambda$. That formula (1) yields a FS is seen by direct verification:

$$
\begin{aligned}
P(\partial) E & =\frac{1}{\overline{P_{m}(\eta)}} \int_{\lambda \in \mathbb{C},|\lambda|=1} \lambda^{m} P(\partial)\left(\mathrm{e}^{\lambda \eta x} \mathcal{F}^{-1}\left(\frac{\overline{P(\mathrm{i} \xi+\lambda \eta)}}{P(\mathrm{i} \xi+\lambda \eta)}\right)\right) \frac{\mathrm{d} \lambda}{2 \pi \mathrm{i} \lambda} \\
& =\frac{1}{\overline{P_{m}(\eta)}} \int_{\lambda \in \mathbb{C},|\lambda|=1} \lambda^{m} \mathrm{e}^{\lambda \eta x}\left(P(\partial+\lambda \eta) \mathcal{F}^{-1}\left(\frac{\overline{P(\mathrm{i} \xi+\lambda \eta)}}{P(\mathrm{i} \xi+\lambda \eta)}\right)\right) \frac{\mathrm{d} \lambda}{2 \pi \mathrm{i} \lambda} \\
& =\frac{1}{\overline{P_{m}(\eta)}} \int_{\lambda \in \mathbb{C},|\lambda|=1} \lambda^{m} \mathrm{e}^{\lambda \eta x} \mathcal{F}^{-1}(\overline{P(\mathrm{i} \xi+\lambda \eta)}) \frac{\mathrm{d} \lambda}{2 \pi \mathrm{i} \lambda} \\
& =\frac{1}{\overline{P_{m}(\eta)}} \int_{\lambda \in \mathbb{C},|\lambda|=1} \lambda^{m} \mathrm{e}^{\lambda \eta x} \overline{P(\partial+\lambda \eta)} \delta \frac{\mathrm{d} \lambda}{2 \pi \mathrm{i} \lambda}
\end{aligned}
$$

$$
=\frac{1}{\overline{P_{m}(\eta)}} \int_{\lambda \in \mathbb{C},|\lambda|=1} \lambda^{m} \mathrm{e}^{\lambda \eta x}\left[\overline{\lambda^{m}} \overline{P_{m}(\eta)} \delta+\sum_{k=0}^{m-1} \overline{\lambda^{k}} Q_{k}(\partial) \delta\right] \frac{\mathrm{d} \lambda}{2 \pi \mathrm{i} \lambda}=\delta
$$

## 2. Duality and microlocal analysis

An important step in the calculation of FSs consists in the determination of its singular support. Here I would like to sketch a connexion of microlocal analysis with Plücker's theory of dual algebraic curves. This relation is also at the heart of the Atiyah-Bott-Gårding theory, but applies to non-hyperbolic operators of principal type as well. For operators with, e.g., double characteristics, things are more difficult, see the articles by Tsuji Mikio [50,51], and by Lars Hörmander [19].

Let $P(\partial)$ be a real homogeneous operator of the degree $m$ and of principal type, i.e., $\forall \xi \in \mathbf{R}^{n} \backslash\{0\}: \mathrm{d} P(\xi) \neq 0$. Then we can easily solve the division problem $P(\omega) \cdot \Phi=1$ on the sphere $\mathbf{S}^{n-1}$ by $\Phi=\operatorname{vp} \frac{1}{P(\omega)} \in \mathcal{D}^{\prime}\left(\mathbf{S}^{n-1}\right)$, which is defined through

$$
\langle\psi, \Phi\rangle=\lim _{\varepsilon \searrow 0} \int_{|P(\omega)|>\varepsilon} \frac{\psi(\omega)}{P(\omega)} \mathrm{d} \sigma(\omega), \quad \psi \in \mathcal{D}\left(\mathbf{S}^{n-1}\right)
$$

From this we obtain $T \in \mathcal{S}^{\prime}\left(\mathbf{R}^{n}\right)$ with $P(\xi) \cdot T=1$ by putting $T=\operatorname{Pf}_{\lambda=-m}\left[\Phi(\xi /|\xi|)|\xi|^{\lambda}\right]$. The distribution $T$ is homogeneous in $\mathbf{R}^{n} \backslash\{0\}$ and we can describe its wave front set WF $T$ quite explicitly: Near a zero $\xi_{0} \in \mathbf{R}^{n} \backslash\{0\}$ of $P$, we use $y_{1}=P(\xi)$ as a coordinate and obtain, since WF $T$ is defined intrinsically in the cotangent space $T^{*} \mathbf{R}^{n}$, that

$$
\begin{aligned}
\mathrm{WF} T \cap T^{*}\left(\mathbf{R}^{n} \backslash\{0\}\right)=\{(\xi, x) ; & \xi \in \mathbf{R}^{n} \backslash\{0\}, P(\xi)=0, x=t \cdot \mathrm{~d} P(\xi), \\
& t \in \mathbf{R} \backslash\{0\}\} .
\end{aligned}
$$

Making use of the following theorem, which goes back to the school of Sato, we obtain a precise description of WF $E$, where $E:=\left(\mathrm{i}^{m} /(2 \pi)^{n}\right) \mathcal{F} T$ is a FS of $P(\partial)$.

Theorem [18, Theorem 8.1.8]. Let $u \in \mathcal{D}^{\prime}\left(\mathbf{R}^{n}\right)$ be homogeneous in $\mathbf{R}^{n} \backslash\{0\}$ and identify $T^{*} \mathbf{R}^{n}$ with $\mathbf{R}^{2 n}$. Then

$$
\begin{aligned}
&(x, \xi) \in \operatorname{WF}(u) \Leftrightarrow(\xi,-x) \in \operatorname{WF}(\mathcal{F} u) \quad \text { if } \xi \neq 0, x \neq 0, \\
& x \in \operatorname{supp} u \Leftrightarrow(0,-x) \in \operatorname{WF}(\mathcal{F} u) \quad \text { if } x \neq 0, \\
& \xi \in \operatorname{supp} \mathcal{F} u \Leftrightarrow \\
&(0, \xi) \in \operatorname{WF}(u) \quad \text { if } \xi \neq 0 .
\end{aligned}
$$

Hence we conclude that
$\mathrm{WF} E=\left\{(x, \xi) ; \xi \in \mathbf{R}^{n} \backslash\{0\}, x=0\right.$ or $\left.[P(\xi)=0, x=t \cdot \nabla P(\xi), t \in \mathbf{R} \backslash\{0\}]\right\}$, where $\nabla P=\left(\frac{\partial P}{\partial x_{1}}, \ldots, \frac{\partial P}{\partial x_{n}}\right)^{\mathrm{T}}$. In particular,
$\operatorname{sing} \operatorname{supp} E=\left\{t \cdot \nabla P(\xi) ; t \in \mathbf{R}, \xi \in \mathbf{R}^{n}, P(\xi)=0\right\}$.
This means that $\operatorname{sing} \operatorname{supp} E$ is the algebraic variety dual to the zero variety of $P$. Let us recapitulate this concept from algebraic geometry.

If $V$ is a finite-dimensional vector space over $\mathbf{K}=\mathbf{R}$ or $\mathbf{C}$, then the corresponding projective space is the set of all one-dimensional subspaces in $V$, i.e.,

$$
\mathbf{P}(V)=\{[v] ; v \in V \backslash\{0\}\}, \quad[v]=\mathbf{K} \cdot v .
$$

The projective space $\mathbf{P}\left(V^{*}\right)$ is canonically identified with the set of all subspaces of $V$ of codimension one and is called the dual projective space. If $X \subset \mathbf{P}(V)$ is a hypersurface given as the zero-set of a homogeneous polynomial $P$ as above, i.e., $X=\{[v] \in \mathbf{P}(V)$; $P(v)=0\}$, then the set of tangent planes to $X$ is an algebraic variety in $\mathbf{P}\left(V^{*}\right)$, called the dual hypersurface $X^{*}$. We consider these varieties over $\mathbf{K}=\mathbf{R}$ or $\mathbf{C}$, and denote them by $X$ or $X^{c}, X^{*}$ or $X^{c *}$, respectively. From the above discussion, we obtain in our case $X^{*}=\{[x] ; x \in \operatorname{sing} \operatorname{supp} E \backslash\{0\}\}$.

Trivial example: The cubic $t=s^{3}$ (written projectively as $x_{1}^{3}-x_{2} x_{3}^{2}=0$ with $s=x_{1} / x_{3}$, $t=x_{2} / x_{3}$ ) has a flex at $s=t=0$. The dual curve is, by definition, the collection of all tangent lines, i.e., $t=3 s_{0}^{2}\left(s-s_{0}\right)+s_{0}^{3}=k s+d$ and hence is parametrized by $k=3 s_{0}^{2}$, $d=-2 s_{0}^{3}$. This is Neill's parabola, which has a cusp at the point $k=d=0$ corresponding to $s=t=0$.

In general, flexes and cusps correspond to one another by duality in the case of plane curves. If $\kappa, \delta, b, f$ and $\kappa^{*}, \delta^{*}, b^{*}, f^{*}$ denote the number of cusps, (ordinary) double points, bitangents, flexes of a plane algebraic curve and of its dual, respectively, then the classical Plücker formulae say (cf. [10, p. 280])

$$
\begin{aligned}
& b=\delta^{*}, \quad b^{*}=\delta, \quad f=\kappa^{*}, \quad f^{*}=\kappa, \quad d^{*}=d(d-1)-2 \delta-3 \kappa, \\
& g=\binom{d-1}{2}-\delta-\kappa .
\end{aligned}
$$

Here $d, d^{*}$ are the degrees of our curves and $g$ denotes the genus.
If $X^{c}$ is given by $P\left(\zeta_{1}, \zeta_{2}, \zeta_{3}\right)=0$, we obtain $X^{c *}$ as the set of those projective points $[z] \in \mathbf{P}\left(\mathbf{C}^{3}\right)$ where the two equations $\zeta \cdot z=0, P(\zeta)=0$ have a multiple projective solution [ $\zeta$ ], and thus from the zero set of the discriminant of $P\left(u,-\left(u z_{1}+z_{3}\right) / z_{2}, 1\right)$ with respect to $u$.

## 3. Homogeneous cubic and quartic operators in 3D

As mentioned earlier, I. Fredholm calculated the FS of $\partial_{1}^{4}+\partial_{2}^{4}+\partial_{3}^{4}$, whereas N. Zeilon failed to find an explicit representation for a FS of $\partial_{1}^{3}+\partial_{2}^{3}+\partial_{3}^{3}$. In later years, Herglotz, Petrovsky, Garnir, etc., explicitly calculated FSs for products of wave and Laplace operators, but, up to $1997, \partial_{1}^{4}+\partial_{2}^{4}+\partial_{3}^{4}$ remained the only irreducible homogeneous operator of degree $>2$, the FS of which was known. In 1997, I succeeded in representing a FS of $\partial_{1}^{3}+\partial_{2}^{3}+\partial_{3}^{3}$ (which I called "Zeilon's operator") by elliptic integrals, and in 1998, I generalized the result to operators of the form $\partial_{1}^{3}+\partial_{2}^{3}+\partial_{3}^{3}+3 a \partial_{1} \partial_{2} \partial_{3}, a \in \mathbf{R} \backslash\{-1\}$ (cf. [53,54]). Let me describe the main result.

According to Newton's classification of real elliptic curves, the non-singular real homogeneous polynomials $P(\xi)$ of third order in three variables are divided into two types according to whether the real projective curve $\left\{[\xi] \in \mathbf{P}\left(\mathbf{R}^{3}\right): P(\xi)=0\right\}$ consists of one or
of two connected components, respectively. In Hesse's normal form, all non-singular real cubic curves are-up to linear transformations-given by

$$
P_{a}(\xi)=\xi_{1}^{3}+\xi_{2}^{3}+\xi_{3}^{3}+3 a \xi_{1} \xi_{2} \xi_{3}, \quad a \in \mathbf{R} \backslash\{-1\}
$$

(Intuitively, this comes from the fact that a homogeneous cubic polynomial in 3 variables, i.e., $P(\xi)=\sum_{\alpha \in \mathbf{N}_{0}^{3},|\alpha|=3} c_{\alpha} \xi^{\alpha}$ has $\binom{3+2}{3}=10$ coefficients and $\operatorname{dimgl}\left(\mathbf{R}^{3}\right)=9$ and hence the (Teichmüller) space of elliptic curves is one-dimensional.) Let $X_{a}:=$ $\left\{[\xi] \in \mathbf{P}\left(\mathbf{R}^{3}\right) ; P_{a}(\xi)=0\right\}$ denote the real projective variety defined by $P_{a}$. For $a>-1$, $X_{a}$ is connected, whereas, for $a<-1, X_{a}$ consists of two components (cf. Fig. 1). The corresponding operators $P_{a}(\partial)$ also differ from the physical viewpoint: For $a<-1$, every projective line through $[1,1,1]$ intersects $X_{a}$ in three different projective points and thus $P_{a}$ is strongly hyperbolic in the direction $(1,1,1)$, for $a>-1, P_{a}$ is not hyperbolic in any direction, nor is it an evolution operator.

We define the fundamental solution $E_{a}$ of $P_{a}(\partial)$ as the Fourier transform of the homogeneous distribution which is of order -3 and has $\mathrm{vp} \frac{1}{P_{a}(\omega)} \in \mathcal{D}^{\prime}\left(\mathbf{S}^{2}\right)$ as its restriction to the sphere. According to Section 2, the (analytic) singular support of $E_{a}$ is the dual curve of $X_{a}$, i.e.,

$$
\operatorname{sing} \operatorname{supp} E_{a}=\operatorname{sing} \operatorname{supp}_{\mathrm{A}} E_{a} \quad \text { and } \quad\left[\operatorname{sing} \operatorname{supp} E_{a} \backslash\{0\}\right]=X_{a}^{*} .
$$

By the classical Plücker formulae, $X_{a}^{*}$ is an algebraic curve of degree 6. Its complexification has nine cusps, three of which are real in correspondence with the three flexes of $X_{a}$ (cf. Fig. 2). Explicitly, we have sing $\operatorname{supp} E_{a}=\left\{x \in \mathbf{R}^{3} ; A_{a}(x)=0\right\}$, where

$$
\begin{align*}
A_{a}(x):= & 3 a\left(a^{3}+4\right) x_{1}^{2} x_{2}^{2} x_{3}^{2}+4\left(a^{3}+1\right)\left(x_{1}^{3} x_{2}^{3}+x_{1}^{3} x_{3}^{3}+x_{2}^{3} x_{3}^{3}\right) \\
& +6 a^{2} x_{1} x_{2} x_{3}\left(x_{1}^{3}+x_{2}^{3}+x_{3}^{3}\right)-\left(x_{1}^{3}+x_{2}^{3}+x_{3}^{3}\right)^{2} . \tag{2}
\end{align*}
$$

If $a<-1$, then $P_{a}$ is hyperbolic with respect to $(1,1,1)$, and $X_{a}^{*}$ consists of two conical surfaces which are the respective duals of the two components of $X_{a}$. Let $F_{a}$ denote the


Fig. 1. $\left\{\left(\xi_{1}, \xi_{2}\right):\left[\xi_{1}, \xi_{2}, 1\right] \in X_{a}\right\}$ for $a=-2$ and for $a=0$.


Fig. 2. $\left\{x \in \mathbf{R}^{3} ;[x] \in X_{a}^{*}, x_{1}+x_{2}+x_{3}=1\right\}$ for $a=-10$ and for $a=1 / 3$.
unique fundamental solution of $P_{a}(\partial)$ with support in $\left\{x \in \mathbf{R}^{3}: x_{1}+x_{2}+x_{3} \geqslant 0\right\}$. Then $E_{a}=\frac{1}{2}\left(F_{a}-\check{F}_{a}\right)$, where the superscript ${ }^{2}$ indicates reflection with respect to the origin. Further, we denote by $K_{a}$ the propagation cone of $P_{a}$ with respect to ( $1,1,1$ ), i.e.,

$$
\begin{equation*}
K_{a}:=\text { dual cone of the component of }(1,1,1) \text { in }\left\{x \in \mathbf{R}^{3} ; P_{a}(x) \neq 0\right\} \tag{3}
\end{equation*}
$$

From the Herglotz-Petrovsky-Leray formula, we infer that $F_{a}$ has a Petrovsky lacuna inside the cone

$$
\begin{equation*}
L_{a}:=\left\{x \in K_{a} ; A_{a}(x)>0\right\} \quad(a<-1) . \tag{4.1}
\end{equation*}
$$

Hence sing supp $F_{a}$ consists of $\partial K_{a}$ and of $\partial L_{a}$, which bound a convex and a non-convex cone, respectively (cf. Fig. 2).

If $a>-1$, then still $E_{a}$ has lacunas inside $L_{a}$ and $-L_{a}$, where now we define

$$
\begin{equation*}
L_{a}:=\text { component of }(1,1,1) \text { in }\left\{x \in \mathbf{R}^{3} ; A_{a}(x)>0\right\} \quad(a>-1) . \tag{4.2}
\end{equation*}
$$

In both cases, the fundamental solutions $E_{a}$ are constant inside $L_{a}$ and $-L_{a}$, and we represent these constant values as complete elliptic integrals of the first kind. Moreover, $E_{a}$ is continuous outside the origin.

Outside the lacunas, $E_{a}(x)$ can be represented by elliptic integrals of the first kind. The final result is contained in the following theorem (cf. [53, p. 286]).

Theorem. Let $a \in \mathbf{R} \backslash\{-1\}$. The limit

$$
T_{a}:=\lim _{\varepsilon \searrow 0} \frac{Y\left(\left|\xi_{1}^{3}+\xi_{2}^{3}+\xi_{3}^{3}+3 a \xi_{1} \xi_{2} \xi_{3}\right|-\varepsilon\right)}{\xi_{1}^{3}+\xi_{2}^{3}+\xi_{3}^{3}+3 a \xi_{1} \xi_{2} \xi_{3}}
$$

defines a distribution in $\mathcal{S}^{\prime}\left(\mathbf{R}^{3}\right)$. If $E_{a}:=(\mathrm{i} /(2 \pi))^{3} \mathcal{F} T_{a}$, and $A_{a}, L_{a}$, and, for $a<-1, K_{a}$ are as in (2) (4.1), (4.2), (3), respectively, then
(a) $E_{a}$ is a fundamental solution of $\partial_{1}^{3}+\partial_{2}^{3}+\partial_{3}^{3}+3 a \partial_{1} \partial_{2} \partial_{3}$;
(b) $E_{a}$ is homogeneous of degree 0;
(c) $E_{a}$ is odd and invariant under permutations of the co-ordinates;
(d) $\operatorname{sing} \operatorname{supp} E_{a}=\operatorname{sing} \operatorname{supp}_{\mathrm{A}} E_{a}=\left\{x \in \mathbf{R}^{3} ; A_{a}(x)=0\right\}$;
(e) $E_{a}$ is continuous in $\mathbf{R}^{3} \backslash\{0\}$;
(f) if $a<-1$, then $E_{a}=\frac{1}{2}\left(F_{a}-\check{F}_{a}\right), P_{a}(\partial) F_{a}=\delta$, $\operatorname{supp} F_{a}=K_{a}$;
(g) $E_{a}$ is constant in $L_{a}$ and in $-L_{a}$, and the values $\left.E_{a}\right|_{L_{a}}$ are given by the following complete elliptic integrals of the first kind:

$$
\left.E_{a}\right|_{L_{a}}=-\frac{1}{4 \sqrt{3} \pi} \begin{cases}\int_{\rho}^{\infty} \frac{\mathrm{d} u}{\sqrt{p_{a}(u)}}, & a>-1 \\ \int_{-\infty}^{\rho} \frac{2 \mathrm{~d} u}{\sqrt{p_{a}(u)}}, & a<-1\end{cases}
$$

where $p_{a}(u):=4\left(a^{3}+1\right) u^{3}+9 a^{2} u^{2}+6 a u+1$ and $\rho$ is the smallest real root of $p_{a}(u)$;
(h) let $x \in U_{a}$, where $U_{a}:=\mathbf{R}^{3} \backslash\left(\bar{L}_{a} \cup-\bar{L}_{a}\right)$ if $a>-1$ and $U_{a}:=\stackrel{\circ}{K_{a}} \backslash L_{a}$ if $a<-1$, and denote by $z(x)$ the only simple real root or, if $x$ belongs to one of the co-ordinate axes, the triple root 0 , respectively, of the cubic equation

$$
\begin{aligned}
Q_{a}(x, z):= & A_{a}(x) z^{3}+9\left(a x_{1}^{2}+x_{2} x_{3}\right)\left(a x_{2}^{2}+x_{1} x_{3}\right)\left(a x_{3}^{2}+x_{1} x_{2}\right) z^{2} \\
& +\left[9 a^{2} x_{1}^{2} x_{2}^{2} x_{3}^{2}+6 a\left(x_{1}^{3} x_{2}^{3}+x_{1}^{3} x_{3}^{3}+x_{2}^{3} x_{3}^{3}\right)\right. \\
& \left.+3 x_{1} x_{2} x_{3}\left(x_{1}^{3}+x_{2}^{3}+x_{3}^{3}\right)\right] z \\
& +3 a x_{1}^{2} x_{2}^{2} x_{3}^{2}+x_{1}^{3} x_{2}^{3}+x_{1}^{3} x_{3}^{3}+x_{2}^{3} x_{3}^{3}=0 .
\end{aligned}
$$

Then $z$ is a real-analytic function in $U_{a}$, and

$$
E_{a}(x)=\left.\frac{Y(-1-a)}{2} E_{a}\right|_{L_{a}}+\frac{\operatorname{sign}\left(\widetilde{P}_{a}(x)\right)}{4 \sqrt{3} \pi} \int_{\rho}^{z(x)} \frac{\mathrm{d} u}{\sqrt{p_{a}(u)}}
$$

where $\widetilde{P}_{a}(x):=3\left[\left(a^{3}-2\right) \rho+a^{2}\right] x_{1} x_{2} x_{3}-(3 a \rho+1)\left(x_{1}^{3}+x_{2}^{3}+x_{3}^{3}\right)$.
Sketch of the proof. Applying the residue theorem in the Herglotz-Petrovsky-Leray formula and using some substitution yields

$$
\begin{equation*}
E_{a}(x)=C_{1}+C_{2} \cdot \operatorname{Im} \int_{\gamma(x)} \Omega, \quad x \in U_{a} \tag{5}
\end{equation*}
$$

where $C_{1}, C_{2}$ are constants, $\gamma(x)$ is a path in the elliptic curve $X_{a}^{c}:=\left\{[\zeta] \in \mathbf{P}\left(\mathbf{C}^{3}\right)\right.$; $\left.P_{a}(\zeta)=0\right\}$ starting at some fixed point and leading to $[y(x)] \in X_{a}^{c}$ defined by $x \cdot y(x)=0$ and $\operatorname{Im} y_{1}(x)>0$, say. Furthermore, $\Omega$ is a generator of the space of holomorphic oneforms on $X_{a}^{c}$. Then the addition theorem for elliptic functions (respectively Abel's theorem for elliptic curves) is applied.

Remarks. Interestingly, it follows from this theorem that the level surfaces of $E_{a}$ are algebraic. Up to present, there is no theoretical explanation for this fact.

In the papers [56,57], we deduce similar formulae for elliptic, respectively, hyperbolic quartic operators of the form $P(\partial)=\sum_{j, k=1}^{3} c_{j k} \partial_{j}^{2} \partial_{k}^{2}$. A typical picture of the slowness surface $X$ and of the dual surface $X^{*}$ for such an operator is given in Fig. 3.


Fig. 3. $X$ and $X^{*}$ for $\partial_{1}^{4}+\partial_{2}^{4}+\partial_{3}^{4}+2 a \partial_{1}^{2} \partial_{2}^{2}+2 b \partial_{3}^{2}\left(\partial_{1}^{2}+\partial_{2}^{2}\right)$ with $a=-0.7$ and $b=-1.2$.

In the regions $A, B$, the fundamental solution $E$ is given by incomplete elliptic integrals of the first kind, in the Petrovsky lacunas $L$ it is given by linear functions the coefficients of which are complete elliptic integrals of the first kind. We refer to [57] for details.

Note that the Riemann surfaces defined by $P(z)=0$ in this case have genus 3 , but $E$ is still given by elliptic integrals. This comes from the fact that $E$ is represented by sums of Abelian integrals in analogy with the imaginary part appearing in formula (5) above.

## 4. The system of crystal optics

If $\mathcal{H}$ denotes the magnetic field, and $\mathcal{J}$ denotes the density of current, and

$$
\varepsilon=\left(\begin{array}{ccc}
\varepsilon_{1} & 0 & 0 \\
0 & \varepsilon_{2} & 0 \\
0 & 0 & \varepsilon_{3}
\end{array}\right)
$$

$\mu, c, \varepsilon_{j}$ being positive constants, then

$$
\left(I_{3} \partial_{t}^{2}+A(\nabla)\right) \mathcal{H}=\frac{4 \pi c}{\mu} \operatorname{rot}\left(\varepsilon^{-1} \mathcal{J}\right)
$$

with the symmetric matrix

$$
A(\xi)=\left(\begin{array}{ccc}
-d_{3} \xi_{2}^{2}-d_{2} \xi_{3}^{2} & d_{3} \xi_{1} \xi_{2} & d_{2} \xi_{1} \xi_{3} \\
d_{3} \xi_{1} \xi_{2} & -d_{3} \xi_{1}^{2}-d_{1} \xi_{3}^{2} & d_{1} \xi_{2} \xi_{3} \\
d_{2} \xi_{1} \xi_{3} & d_{1} \xi_{2} \xi_{3} & -d_{1} \xi_{2}^{2}-d_{2} \xi_{1}^{2}
\end{array}\right)
$$

where we have set $d_{j}=c^{2} /\left(\mu \varepsilon_{j}\right), j=1,2,3$. Particularly important for systems of PDOs is the determinant operator. We have

$$
\operatorname{det}\left(I_{3} \partial_{t}^{2}+A(\nabla)\right)=\partial_{t}^{2} R(\partial)
$$

with

$$
R(\tau, \xi)=\tau^{4}-\tau^{2} \sum_{j=1}^{3} \xi_{j}^{2}\left(d_{j+1}+d_{j+2}\right)+|\xi|^{2} \sum_{j=1}^{3} \xi_{j}^{2} d_{j+1} d_{j+2}
$$

(where we define $d_{4}=d_{1}, d_{5}=d_{2}$ ). The slowness surface

$$
X=\left\{[(\tau, \xi)] \in \mathbf{P}\left(\mathbf{R}^{4}\right) ; R(\tau, \xi)=0\right\} \cong\left\{\xi \in \mathbf{R}^{3} ; R(1, \xi)=0\right\}
$$

is called "Fresnel's surface." If the positive constants $d_{1}, d_{2}, d_{3}$ are pairwise different, then $X$ is homeomorphic to two disjoint spheres glued together at four points, which two by two are pairwise opposite and span the "optical axes." In this case, $R$ is an irreducible polynomial. If two of the $d_{j}$ are equal, i.e., if $d_{1}=d_{2}=1$ and $d_{3}=d \neq 1$ without loss of generality, then the crystal is called "uniaxial," because in this case the optical axes coincide. Then $X$ is made up of a sphere and of an ellipsoid touching each other at the two points on the optical axis.

Use of the matrix version of the Herglotz-Petrovsky formula yields a representation of $E$ by Abelian integrals (cf. [39, 2.2.2, p. 327], [26, p. 3318])

$$
\begin{equation*}
E(t, x)=-\frac{Y(t)}{4 \pi^{2}} \partial_{t} \int_{C_{t, x}} \frac{P(1, \xi)^{\mathrm{ad}} \operatorname{sign}\left(\left(\partial_{\tau} \operatorname{det} P\right)(1, \xi)\right)}{\left|x_{3}\left(\partial_{2} \operatorname{det} P\right)(1, \xi)-x_{2}\left(\partial_{3} \operatorname{det} P\right)(1, \xi)\right|}\left|\mathrm{d} \xi_{1}\right| \tag{6}
\end{equation*}
$$

where

$$
C_{t, x}:=\left\{\xi \in \mathbf{R}^{3} ; \operatorname{det} P(1, \xi)=0, t+\xi^{\mathrm{T}} \cdot x=0\right\} \quad\left(\text { for }(t, x) \in \mathbf{R}^{4}\right) .
$$

Unfortunately, in the case of crystal optics (i.e., $P(\partial)=I_{3} \partial_{t}^{2}+A(\nabla), A$ as above), an explicit evaluation (in terms of higher transcendental functions) of formula (6) has not yet been achieved. Let me describe what is known so far [39, 3.4 and 4.3]:

If $K$ denotes the support of $E$, then $K$ is the dual cone of the connectivity component of $(1,0)$ in $\left\{(\tau, \xi) \in \mathbf{R}^{4} ; R(\tau, \xi) \neq 0\right\}$. The singular support of $E$ consists of the four "Hamiltonian circles" and of $X^{*}$, where $X^{*}$ is the dual surface to the slowness surface $X$. It turns out that

$$
\begin{aligned}
& X^{*}=\left\{[(t, x)] \in \mathbf{P}\left(\mathbf{R}^{4}\right) ;\right. \\
&\left.t^{4}-t^{2} \sum_{j=1}^{3} x_{j}^{2}\left(d_{j+1}^{-1}+d_{j+2}^{-1}\right)+|x|^{2} \sum_{j=1}^{3} x_{j}^{2} /\left(d_{j+1} d_{j+2}\right)=0\right\}
\end{aligned}
$$

and hence $X^{*}$ is given by an equation analogous to that of Fresnel's surface $X$. The intersection of $X^{*}$ with a plane through the optical axes consists of a circle intersecting an ellipse, see Fig. 4.

The fundamental matrix $E$ is explicitly known in the inner region $J$. There

$$
E=\operatorname{vp} \frac{t Y(t)}{4 \pi|x|^{3}}\left(I_{3}-\frac{3 x \cdot x^{\mathrm{T}}}{|x|^{2}}\right)+\frac{1}{3}(t Y(t) \otimes \delta(x)) I_{3}
$$

Furthermore, one can calculate the delta terms in $E$, and $E$ is known in the uniaxial case. In this case, the circle and the ellipse in Fig. 4 touch each other and $E$ can be expressed


Fig. 4. Section of $\operatorname{sing} \operatorname{supp} E$ through the optical axes.
by delta terms and algebraic functions [39, Proposition 3, p. 342]. For the biaxial case, however, $E$ is given in $K \backslash J$ by Abelian integrals over curves of genus 3, the moduli of which depend on $(t, x)$. Up to now, there is no representation by higher transcendental functions known.

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[^0]:    E-mail address: wagner@mat1.uibk.ac.at.
    URL: http://techmath.uibk.ac.at/mathematik/wagner.

