# On the Fourier transform of rotationally invariant distributions 

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#### Abstract

We present an extension of the Poisson-Bochner formula for the Fourier transform of rotationally invariant distributions by analytic continuation "with respect to the dimension". As application of this extension, a new derivation of the fundamental solution of the Euler-Poisson-Darboux operator is given.


Keywords Fourier transform • Poisson-Bochner formula • Euler-Poisson-Darboux operator • Analytic continuation

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## 1 Introduction and notation

If $f \in L^{1}\left(\mathbb{R}^{n}\right)$ is rotationally invariant, i.e., if $f(x)=g(|x|)$ with $g(\rho) \rho^{n-1} \in L^{1}((0, \infty))$, then the classical Poisson-Bochner formula expresses the Fourier transform $\mathcal{F} f \in \mathcal{C}_{0}\left(\mathbb{R}^{n}\right)$ of $f$ by the integral

$$
\begin{equation*}
(\mathcal{F} f)(\xi)=(2 \pi)^{n / 2}|\xi|^{-n / 2+1} \int_{0}^{\infty} g(\rho) \rho^{n / 2} J_{n / 2-1}(\rho|\xi|) \mathrm{d} \rho, \quad \xi \neq 0 \tag{1.1}
\end{equation*}
$$

see [36, (VII, 7; 22), p. 259], [37, Thm. 3.10, p. 158]. Our naming "Poisson-Bochner formula" is motivated by the generalization of the formula (1.1) for dimensions $n=2,3$ (discovered by Poisson and Cauchy, see footnote 109 in [5, p. 226]) to dimensions $n \geq 4$ by S. Bochner in [5, Satz 56, p. 186].

A generalization for functions in weighted $L^{1}$-spaces, i.e., for $g \in L_{\mathrm{loc}}^{1}((0, \infty))$ fulfilling

$$
\begin{equation*}
\int_{0}^{\infty}|g(\rho)| \rho^{n-1}(1+\rho)^{(1-n) / 2} \mathrm{~d} \rho<\infty \tag{1.2}
\end{equation*}
$$

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is given in [33, Lemma 25.1, p. 358; Engl. transl. p. 485]. (In fact, if (1.2) is satisfied, then $\mathcal{F} f$ is continuous outside the origin and formula (1.1) is valid for $\xi \neq 0$ by applying Lebesgue's theorem on dominated convergence to $\lim _{N \rightarrow \infty}\langle\phi(\xi), \mathcal{F}(f(x) \cdot Y(N-|x|))\rangle$ for $\phi \in \mathcal{D}\left(\mathbb{R}^{n} \backslash\{0\}\right)$.) A further generalization by means of partial integration can be found in [24,25], see also [31, Ex. 1.6.13 (a), p. 102]. A limit representation of the Fourier transform of radial temperate distributions is given in [26, (2.109), p. 140].

Let us mention that, e.g., the forward fundamental solution $E$ of the wave operator $\partial_{t}^{2}-\Delta_{n}$ is given by

$$
E=(2 \pi)^{-n} Y(t) \mathcal{F}_{x}\left(\frac{\sin (t|x|)}{|x|}\right),
$$

see Examples 2.2, 2.4 below. In this case, for $t>0, g(\rho)=(2 \pi)^{-n} \rho^{-1} \sin (t \rho)$ and $f(x)=g(|x|)$ is not integrable nor does $g$ satisfy condition (1.2). In order to calculate this important Fourier transform, different approximation methods were used, compare [13, pp. 177-183], [38, p. 51], [31, Ex. 1.6.17 (a), p. 106, Ex. 2.3.6, p. 141].

The main purpose of this paper consists in generalizing formula (1.1) so as to yield a representation of the Fourier transform $\mathcal{F} S$ for arbitrary radially symmetric temperate distributions $S$. This is done by analytic continuation with respect to the index $\lambda=\frac{n}{2}-1$ of the Bessel function in (1.1), see Theorems 2.1, 2.3. So in a way, we use "analytic continuation with respect to the dimension $n$ " of the underlying space $\mathbb{R}^{n}$. Heuristically, this procedure goes back, at least, to A. Weinstein, comp. [39, p. 44]: "The viewpoint of spaces of 'fractional dimensions' due to Weinstein is very fruitful and led to fundamental solutions in the large of the iterated EPD-equation."

In [11, p. 8], the Bochner transform $T_{n}$ is defined by

$$
T_{n} \varphi(r)=\frac{2 \pi}{r^{n / 2-1}} \int_{0}^{\infty} J_{n / 2-1}(2 \pi r \rho) \rho^{n / 2} \varphi(\rho) \mathrm{d} \rho
$$

for suitable functions $\varphi$ and $n \in \mathbb{N}$. Whereas in [16] the connection between $T_{n}$ and $T_{n+2}$ is rederived, see [33, (25.14'), Lemma 25.1', p. 359; Engl. transl. p. 486] and [32, p. 270], and in [12,27], the general connection between $T_{n}$ and $T_{n+q}, q \in \mathbb{N}$, is investigated, the present study is concerned with the analytic continuation of the function $\lambda \mapsto T_{\lambda}$ for complex $\lambda$.

In order to illustrate our method, we first apply it to the wave equation (Examples 2.2, 2.4) and then, in Sect. 3, to the Euler-Poisson-Darboux equation. In Propositions 3.2, 3.3, we derive in this way the fundamental solution $E$ of the EPD-operator. (For the concept of fundamental solutions of linear partial differential operators with non-constant coefficients, see [36, pp. 138-142], [23, p. 29], [9, pp. 11-14].) This more complicated fundamental solution was given in [10] and verified therein by series expansion, see $[3,4]$ for a recapitulation. Our deduction of $E$ based on the analytic continuation of the Poisson-Bochner formula is different from that in $[3,4,10$ ] and seems to be new, comp. [2, p. 478]: "We do not know how to obtain an explicit formula (or formulas) for the inverse Fourier transform of $\tilde{F}(\xi, y ; b)$ when $b \neq 0$, a problem that merits to be investigated."

Let us introduce some notation. We employ the standard notation for the distribution spaces $\mathcal{D}^{\prime}, \mathcal{S}^{\prime}, \mathcal{E}^{\prime}$, the dual spaces of the spaces $\mathcal{D}, \mathcal{S}, \mathcal{E}$ of "test functions", of "rapidly decreasing functions" and of $C^{\infty}$ functions, respectively, see [18,20,36]. In order to display the active variable in a distribution, say $x \in \mathbb{R}^{n}$, we use notation as $T(x)$ or $T \in \mathcal{D}^{\prime}\left(\mathbb{R}_{x}^{n}\right)$. Furthermore, we use the spaces $\mathcal{D}_{L^{p}}, \mathcal{D}_{L^{p}}^{\prime}, 1 \leq p \leq \infty, \mathcal{O}_{M}, \mathcal{O}_{C}^{\prime}$, which were introduced in [36, Ch. VI, § 8, p. 199; Ch. VII, § 5, p. 243], and we set $\mathcal{S}_{r}^{\prime}\left(\mathbb{R}^{n}\right)=\left\{S \in \mathcal{S}^{\prime}\left(\mathbb{R}^{n}\right) ; S\right.$ is radially symmetric $\}$.

For the evaluation of a distribution $T \in E^{\prime}$ on a test function $\phi \in E$, we use angle brackets, i.e., $\langle\phi, T\rangle$ or, more precisely ${ }_{E}\langle\phi, T\rangle_{E^{\prime}}$. More generally, if $\phi \in E \hat{\otimes} F$ and $T \in E^{\prime}$
for distribution spaces $E, F$, then ${ }_{E \hat{\otimes} F}\langle\phi, T\rangle_{E^{\prime}}$ symbolizes the vector-valued scalar product $(E \hat{\otimes} F) \times E^{\prime} \rightarrow F$, see $[34,35]$ for more information on vector-valued distributions. (In all tensor products of this study, both factors are complete and at least one of the factors is nuclear and hence $E \hat{\otimes}_{\pi} F=E \hat{\otimes}_{\epsilon} F$ and we simply write $E \hat{\otimes} F$.)

The Heaviside function is denoted by $Y$, see [36, p. 36], and we set

$$
\begin{equation*}
\chi^{\mu}(t)=\frac{Y(t) t^{\mu}}{\Gamma(\mu+1)} \in L_{\mathrm{loc}}^{1}\left(\mathbb{R}_{t}^{1}\right) \text { for } \mu \in \mathbb{C} \text { with } \operatorname{Re} \mu>-1 . \tag{1.3}
\end{equation*}
$$

The function $\mu \mapsto \chi^{\mu}$ can be analytically continued in $\mathcal{S}^{\prime}\left(\mathbb{R}^{1}\right)$ and thus yields an entire function

$$
\chi: \mathbb{C} \longrightarrow \mathcal{S}^{\prime}\left(\mathbb{R}^{1}\right): \mu \longmapsto \chi^{\mu}
$$

see $\left[18\right.$, (3.2.17), p. 73]. We write $\delta_{\tau}(t) \in \mathcal{D}^{\prime}\left(\mathbb{R}_{t}^{1}\right), \tau \in \mathbb{R}$, for the delta distribution with support in $\tau$, which is the derivative of $Y(t-\tau)$, i.e., $\left\langle\phi, \delta_{\tau}\right\rangle=\phi(\tau)$ for $\phi \in \mathcal{D}\left(\mathbb{R}^{1}\right)$. In contrast, $\delta \in \mathcal{D}^{\prime}\left(\mathbb{R}^{n}\right)$ without any subscript stands for the delta distribution at the origin.

The pull-back $h^{*} T=T \circ h \in \mathcal{D}^{\prime}(\Omega)$ of a distribution $T$ in one variable $t$ with respect to a submersive $\mathcal{C}^{\infty}$ function $h: \Omega \rightarrow \mathbb{R}, \Omega \subset \mathbb{R}^{n}$ open, is defined as in [14, (7.2.4/5), p. 82] or in [31, Def. 1.2.12, p. 19], i.e.,

$$
\begin{equation*}
\left\langle\phi, h^{*} T\right\rangle=\left\langle\frac{\mathrm{d}}{\mathrm{~d} t}\left(\int_{\Omega} Y(t-h(x)) \phi(x) \mathrm{d} x\right), T\right\rangle, \quad \phi \in \mathcal{D}(\Omega) . \tag{1.4}
\end{equation*}
$$

We use the Fourier transform $\mathcal{F}$ in the form

$$
(\mathcal{F} \phi)(\xi):=\int \mathrm{e}^{-\mathrm{i} \xi x} \phi(x) \mathrm{d} x, \quad \phi \in \mathcal{S}\left(\mathbb{R}^{n}\right)
$$

this being extended to $\mathcal{S}^{\prime}$ by continuity. We write $\left|\mathbb{S}^{n-1}\right|$ for the hypersurface area $2 \pi^{n / 2} / \Gamma\left(\frac{n}{2}\right)$ of the unit sphere $\mathbb{S}^{n-1}$ in $\mathbb{R}^{n}$. For $j \in \mathbb{N}$ and $w \in \mathbb{C}$, we use Pochhammer's symbol $(w)_{0}=1$, $(w)_{j}=w \cdot(w+1) \cdots \cdot(w+j-1) . J_{\lambda}$ and $N_{\lambda}$ denote, as usual, the Bessel functions of the first and of the second kind.

## 2 Analytic continuation of the Poisson-Bochner formula

Let us first rewrite (1.1) in a more symmetrical fashion by the following $n$-dimensional integral, still under the assumptions that $f \in L^{1}\left(\mathbb{R}^{n}\right)$ and $f$ is radially symmetric:

$$
\begin{equation*}
(\mathcal{F} f)(\xi)=2^{n / 2-1} \Gamma\left(\frac{n}{2}\right) \int_{\mathbb{R}^{n}} f(x) \frac{J_{n / 2-1}(|x| \cdot|\xi|)}{(|x| \cdot|\xi|)^{n / 2-1}} \mathrm{~d} x \tag{2.1}
\end{equation*}
$$

We note incidentally that formula (2.1) allows a generalization (which follows, e.g., by density) for $S \in \mathcal{D}_{L^{1}}^{\prime}\left(\mathbb{R}^{n}\right) \cap \mathcal{S}_{r}^{\prime}\left(\mathbb{R}^{n}\right)$, i.e., for radially symmetric integrable distributions $S$. Then $\mathcal{F} S$ is a continuous function given by

$$
\begin{equation*}
(\mathcal{F} S)(\xi)=2^{n / 2-1} \Gamma\left(\frac{n}{2}\right) \cdot \mathcal{D}_{L^{1}}^{\prime}\left(\mathbb{R}_{x}^{n}\right)\left\langle S(x), \frac{J_{n / 2-1}(|x| \cdot|\xi|)}{(|x| \cdot|\xi|)^{n / 2-1}}\right\rangle_{\mathcal{D}_{L^{\infty}\left(\mathbb{R}_{x}^{n}\right)}}, \quad \xi \in \mathbb{R}^{n} \tag{2.2}
\end{equation*}
$$

As can be derived from [17, p. 538], the kernel

$$
K(x, \xi)=\frac{J_{n / 2-1}(|x| \cdot|\xi|)}{(|x| \cdot|\xi|)^{n / 2-1}} \in \mathcal{O}_{M}\left(\mathbb{R}^{2 n}\right)
$$

belongs to the completed tensor product $\mathcal{S}\left(\mathbb{R}_{x}^{n}\right) \hat{\otimes} \mathcal{S}^{\prime}\left(\mathbb{R}_{\xi}^{n}\right)$, and therefore the Fourier transform of $S \in \mathcal{S}_{r}^{\prime}\left(\mathbb{R}^{n}\right)$ can be written in the form

$$
\begin{equation*}
(\mathcal{F} S)(\xi)=2^{n / 2-1} \Gamma\left(\frac{n}{2}\right) \cdot \mathcal{S}^{\prime}\left(\mathbb{R}_{x}^{n}\right)\langle S(x), K(x, \xi)\rangle_{\mathcal{S}\left(\mathbb{R}_{x}^{n}\right) \hat{\otimes} \mathcal{S}^{\prime}\left(\mathbb{R}_{\xi}^{n}\right)} \tag{2.3}
\end{equation*}
$$

(by applying [35, Prop. 4, p. 41]).
Note that formula (2.3) allows to represent $\mathcal{F}_{x}\left(|x|^{-1} \sin (t|x|)\right)$ by the $\mathcal{S}^{\prime}$-valued scalar product

$$
\begin{equation*}
\mathcal{F}_{x}\left(\frac{\sin (t|x|)}{|x|}\right)(\xi)=2^{n / 2-1} \Gamma\left(\frac{n}{2}\right) \cdot \mathcal{S}^{\prime}\left(\mathbb{R}_{x}^{n}\right)\left\langle\frac{\sin (t|x|)}{|x|}, \frac{J_{n / 2-1}(|x| \cdot|\xi|)}{(|x| \cdot|\xi|)^{n / 2-1}}\right\rangle_{\mathcal{S}\left(\mathbb{R}_{x}^{n}\right) \hat{\otimes} \mathcal{S}^{\prime}\left(\mathbb{R}_{\xi}^{n}\right)} . \tag{2.4}
\end{equation*}
$$

However, formula (2.4) cannot be evaluated for fixed $\xi$. In the following two theorems, we shall therefore imbed the kernel $K(x, \xi)$ into an analytic family of kernels $K_{\lambda}(x, \xi)$ such that $\mathcal{F} S, S \in \mathcal{S}_{r}^{\prime}\left(\mathbb{R}^{n}\right)$, can be obtained by analytic continuation with respect to $\lambda$. Let us mention that

$$
\frac{J_{\lambda}(|x| \cdot|\xi|)}{(|x| \cdot|\xi|)^{\lambda}} \in \mathcal{O}_{M}\left(\mathbb{R}^{2 n}\right)
$$

depends holomorphically on $\lambda \in \mathbb{C}$ (see below), but that these kernels do not belong to $\mathcal{S}\left(\mathbb{R}_{x}^{n}\right) \hat{\otimes} \mathcal{S}^{\prime}\left(\mathbb{R}_{\xi}^{n}\right)$ and not even to $\mathcal{D}_{L^{1}}^{\prime}\left(\mathbb{R}_{x}^{n}\right) \hat{\otimes} \mathcal{S}^{\prime}\left(\mathbb{R}_{\xi}^{n}\right)$ for $\lambda \in \mathbb{C} \backslash\left(\frac{n}{2}-\mathbb{N}\right)$. This is the reason for the more complicated choices of $K_{\lambda}$ below.

Theorem 2.1 The kernel

$$
K_{\lambda}(x, \xi)=|\xi|^{2 \lambda-n+2} \cdot \frac{J_{\lambda}(|x| \cdot|\xi|)}{(|x| \cdot|\xi|)^{\lambda}} \in \mathcal{S}^{\prime}\left(\mathbb{R}_{x, \xi}^{2 n}\right)
$$

is an entire function of $\lambda$ with values in $\mathcal{S}\left(\mathbb{R}_{x}^{n}\right) \hat{\otimes} \mathcal{S}^{\prime}\left(\mathbb{R}_{\xi}^{n}\right)$. Furthermore, if $S \in \mathcal{S}_{r}^{\prime}\left(\mathbb{R}^{n}\right)$, then

$$
\mathcal{F} S=2^{n / 2-1} \Gamma\left(\frac{n}{2}\right) \cdot U_{n / 2-1}
$$

where the function

$$
\begin{equation*}
U: \mathbb{C} \longrightarrow \mathcal{S}^{\prime}\left(\mathbb{R}_{\xi}^{n}\right): \lambda \longmapsto U_{\lambda}(\xi)=\mathcal{S}^{\prime}\left(\mathbb{R}_{x}^{n}\right)\left\langle S(x), K_{\lambda}(x, \xi)\right\rangle_{\mathcal{S}\left(\mathbb{R}_{x}^{n}\right) \hat{\otimes} \mathcal{S}^{\prime}\left(\mathbb{R}_{\xi}^{n}\right)} \tag{2.5}
\end{equation*}
$$

is entire.
Proof (a) Let us first show that the mapping

$$
\begin{equation*}
\mathbb{C} \longrightarrow \mathcal{O}_{M}\left(\mathbb{R}_{x, \xi}^{2 n}\right): \lambda \longmapsto \frac{J_{\lambda}(|x| \cdot|\xi|)}{(|x| \cdot|\xi|)^{\lambda}} \tag{2.6}
\end{equation*}
$$

is entire. From [15, 8.411.8] and using analytic continuation, we obtain the representation

$$
\frac{J_{\lambda}(|x| \cdot|\xi|)}{(|x| \cdot|\xi|)^{\lambda}}=\frac{1}{2^{\lambda} \sqrt{\pi} \mathcal{E}\left(\mathbb{R}_{t}^{1}\right)}\left\langle\cos (t|x| \cdot|\xi|),\left.\chi^{-1 / 2+\lambda}\left(1-t^{2}\right)\right|_{\mathcal{E}^{\prime}\left(\mathbb{R}_{t}^{1}\right)}\right.
$$

for each $(x, \xi) \in \mathbb{R}^{2 n}$. (For $\chi^{\mu}$ see Sect. 1, in particular (1.3).) Since

$$
\cos (t|x| \cdot|\xi|) \in \mathcal{O}_{M}\left(\mathbb{R}_{t, x, \xi}^{2 n+1}\right)=\mathcal{O}_{M}\left(\mathbb{R}_{t}^{1}\right) \hat{\otimes} \mathcal{O}_{M}\left(\mathbb{R}_{x, \xi}^{2 n}\right)
$$

see [34, Prop. 28, p. 98], and since

$$
\mathbb{C} \longrightarrow \mathcal{E}^{\prime}\left(\mathbb{R}_{t}^{1}\right): \lambda \longmapsto \chi^{-1 / 2+\lambda}\left(1-t^{2}\right)
$$

is entire and $\mathcal{E}^{\prime} \subset \mathcal{O}_{M}^{\prime}$, we conclude that also the mapping in (2.6) is entire, see [35, Prop. 4, p. 41].
(b) The distribution-valued function

$$
F:\{\lambda \in \mathbb{C} ; \operatorname{Re} \lambda>-1\} \longrightarrow \mathcal{S}^{\prime}\left(\mathbb{R}_{\xi}^{n}\right): \lambda \longmapsto F_{\lambda}=|\xi|^{2 \lambda-n+2} \in L_{\mathrm{loc}}^{1}\left(\mathbb{R}_{\xi}^{n}\right)
$$

can analytically be continued to $\mathbb{C} \backslash(-\mathbb{N})$ and has simple poles for $\lambda=-k, k \in \mathbb{N}$, with the residues

$$
\operatorname{Res}_{\lambda=-k} F_{\lambda}=\frac{\pi^{n / 2} \Delta_{n}^{k-1} \delta}{2^{2(k-1)}(k-1)!\Gamma\left(\frac{n}{2}+k-1\right)},
$$

see [30, Ex. 2.3.1, p. 41]. Therefore, the product

$$
K_{\lambda}(x, \xi)=F_{\lambda}(\xi) \cdot \frac{J_{\lambda}(|x| \cdot|\xi|)}{(|x| \cdot|\xi|)^{\lambda}} \in \mathcal{S}^{\prime}\left(\mathbb{R}_{\xi}^{n}\right) \cdot \mathcal{O}_{M}\left(\mathbb{R}_{x, \xi}^{2 n}\right)
$$

is well-defined and belongs to $\mathcal{O}_{M}\left(\mathbb{R}_{x}^{n}\right) \hat{\otimes} \mathcal{S}^{\prime}\left(\mathbb{R}_{\xi}^{n}\right)$ for $\lambda \in \mathbb{C} \backslash(-\mathbb{N})$, see [34, Props. 20 bis, 28, pp. 70, 98]. Furthermore, $K_{\lambda}(x, \xi)$ is also holomorphic in $\lambda=-k \in-\mathbb{N}$ since

$$
\frac{J_{-k}(|x| \cdot|\xi|)}{(|x| \cdot|\xi|)^{-k}}=(-1)^{k}(|x| \cdot|\xi|)^{k} \cdot J_{k}(|x| \cdot|\xi|)
$$

vanishes of order $2 k$ at $\xi=0$ and hence its product with $\operatorname{Res}_{\lambda=-k} F_{\lambda}(\xi)$ vanishes.
(c) Let us next calculate the partial Fourier transform $\mathcal{F}_{\xi}\left(K_{\lambda}(x, \xi)\right)$. Since the kernel $K_{\lambda}(x, \xi)$ belongs to $\mathcal{O}_{M}\left(\mathbb{R}_{x}^{n}\right) \hat{\otimes} \mathcal{S}^{\prime}\left(\mathbb{R}_{\xi}^{n}\right)$, it is semi-regular in $x$, and it suffices to determine $\mathcal{F}_{\xi}\left(K_{\lambda}(x, \xi)\right)$ for fixed $x \neq 0$. If $x \neq 0$, then $K_{\lambda}(x, \xi)$ is bounded by a constant times $|\xi|^{3 / 2-n+\operatorname{Re} \lambda}$ for $|\xi| \rightarrow \infty$. This implies that condition (1.2) is satisfied upon setting $g(\rho)=$ $|x|^{-\lambda} \rho^{\lambda-n+2} J_{\lambda}(\rho|x|)$ and that $K_{\lambda}(x, \xi) \in L^{1}\left(\mathbb{R}_{\xi}^{n}\right)+L^{2}\left(\mathbb{R}_{\xi}^{n}\right)$ if $n \geq 3$ and $-1<\operatorname{Re} \lambda<$ $\frac{n}{2}-2$. Therefore $\mathcal{F}_{\xi} K_{\lambda}$ belongs to $L_{\text {loc }}^{2}\left(\mathbb{R}_{\xi}^{n}\right)$ for such $x$ and $\lambda$ and the Poisson-Bochner formula applies and represents $\mathcal{F}_{\xi} K_{\lambda}$ by the absolutely convergent integral in (1.1) for $\xi \neq 0$. For $x \neq 0, \xi \neq 0$ and $\lambda$ as above, formula 6.575.1, p. 692, in [15] then yields

$$
\begin{aligned}
\left(\mathcal{F}_{\xi} K_{\lambda}\right)(x, \xi) & =(2 \pi)^{n / 2}|x|^{-\lambda}|\xi|^{-n / 2+1} \int_{0}^{\infty} J_{n / 2-1}(\rho|\xi|) J_{\lambda}(\rho|x|) \rho^{\lambda-(n / 2-2)} \mathrm{d} \rho \\
& =2^{\lambda+2} \pi^{n / 2}|\xi|^{2-n} \chi^{n / 2-\lambda-2}\left(|\xi|^{2}-|x|^{2}\right)
\end{aligned}
$$

(The correct parameter range for formula 6.575.1 in [15] is $-1<\operatorname{Re} \mu<\operatorname{Re}(v+1)$.)
For $x \neq 0$, the function $h(\xi)=|\xi|^{2}-|x|^{2}$ is submersive and hence the composition $\chi^{n / 2-\lambda-2}\left(|\xi|^{2}-|x|^{2}\right) \in \mathcal{S}^{\prime}\left(\mathbb{R}_{\xi}^{n}\right)$ is well-defined, see (1.4), and it is an entire function of $\lambda$. Since $\chi^{n / 2-\lambda-2}\left(|\xi|^{2}-|x|^{2}\right)$ vanishes at $\xi=0$, the product $|\xi|^{2-n} \chi^{n / 2-\lambda-2}\left(|\xi|^{2}-|x|^{2}\right)$ is also well-defined and depends holomorphically on $\lambda$ in $\mathcal{S}^{\prime}\left(\mathbb{R}_{\xi}^{n}\right)$. By analytic continuation, we conclude that $\mathcal{F}_{\xi} K_{\lambda}$ is represented by the continuous function $\mathbb{R}_{x}^{n} \rightarrow \mathcal{S}^{\prime}\left(\mathbb{R}_{\xi}^{n}\right)$ which, for $x \neq 0$, is given by the equation

$$
\begin{equation*}
\mathcal{F}_{\xi} K_{\lambda}=2^{\lambda+2} \pi^{n / 2}|\xi|^{2-n} \chi^{n / 2-\lambda-2}\left(|\xi|^{2}-|x|^{2}\right) \tag{2.7}
\end{equation*}
$$

If $n=2$ or $n=1$, then the same conclusion can be reached by proving (2.7) for $\operatorname{Re} \lambda<-\frac{3}{2}$ with the help of formula (2.2). (Note that $K_{\lambda}(x, \xi) \in \mathcal{D}_{L^{1}}^{\prime}\left(\mathbb{R}_{\xi}^{n}\right)$ for $\operatorname{Re} \lambda<-\frac{3}{2}$ and fixed $x \in \mathbb{R}^{n} \backslash\{0\}$.) Hence (2.7) is valid for $n \in \mathbb{N}, x \neq 0$ and each $\lambda \in \mathbb{C}$.
(d) If $\phi \in \mathcal{S}\left(\mathbb{R}_{\xi}^{n}\right)$, then

$$
\psi(x)=\left\langle\phi(\xi), \mathcal{F}_{\xi} K_{\lambda}\right\rangle=\mathcal{S}\left(\mathbb{R}_{\xi}^{n}\right)\left\langle(\mathcal{F} \phi)(\xi), K_{\lambda}(x, \xi)\right\rangle_{\mathcal{O}_{M}\left(\mathbb{R}_{x}^{n}\right) \hat{\otimes} \mathcal{S}^{\prime}\left(\mathbb{R}_{\xi}^{n}\right)}
$$

belongs to $\mathcal{O}_{M}\left(\mathbb{R}_{x}^{n}\right)$ and (2.7) implies, for $x \neq 0$,

$$
\begin{aligned}
\psi(x) & \left.=\left\langle\phi(\xi), \mathcal{F}_{\xi} K_{\lambda}\right\rangle=\left.2^{\lambda+2} \pi^{n / 2}\langle\phi(\xi),| \xi\right|^{2-n} \chi^{n / 2-\lambda-2}\left(|\xi|^{2}-|x|^{2}\right)\right\rangle \\
& \left.=\left.2^{\lambda+2} \pi^{n / 2}|x|^{n-2 \lambda-2}\langle\phi(|x| \eta),| \eta\right|^{2-n} \chi^{n / 2-\lambda-2}\left(|\eta|^{2}-1\right)\right\rangle .
\end{aligned}
$$

Because the support of the distribution $|\eta|^{2-n} \chi^{n / 2-\lambda-2}\left(|\eta|^{2}-1\right)$ does not contain the origin $\eta=0$, we conclude that $\psi$ is, with all its derivatives, rapidly decreasing for $|x| \rightarrow \infty$ and hence $\psi \in \mathcal{S}\left(\mathbb{R}^{n}\right)$. This means that $\mathcal{F}_{\xi} K_{\lambda}$ and thus also $K_{\lambda}$ belong to $\mathcal{S}\left(\mathbb{R}_{x}^{n}\right) \hat{\otimes} \mathcal{S}^{\prime}\left(\mathbb{R}_{\xi}^{n}\right)$.
(e) We observe that $K_{n / 2-1}(x, \xi)=J_{n / 2-1}(|x| \cdot|\xi|) /(|x| \cdot|\xi|)^{n / 2-1}$ and hence (2.3) implies $\mathcal{F} S=2^{n / 2-1} \Gamma\left(\frac{n}{2}\right) U_{n / 2-1}$. Finally, the map $\lambda \mapsto U_{\lambda}$ in (2.5) is analytic by [35, Prop. 4, p. 41]. This completes the proof.

Example 2.2 Let us illustrate Theorem 2.1 by calculating the forward fundamental solution $E$ of the wave operator $\partial_{t}^{2}-\Delta_{n}$. We consider $E \in \mathcal{S}^{\prime}\left(\mathbb{R}_{t, x}^{n+1}\right)$ as the $\mathcal{C}^{\infty}$ mapping

$$
E:[0, \infty) \longrightarrow \mathcal{S}^{\prime}\left(\mathbb{R}_{x}^{n}\right): t \longmapsto E_{t}(x)
$$

and represent $E_{t}$ by partial Fourier transform, i.e.,

$$
E_{t}=(2 \pi)^{-n} \mathcal{F}_{x}\left(\frac{\sin (t|x|)}{|x|}\right), \quad t \geq 0
$$

compare [31, Ex. 1.6.17, p. 106].
(a) The distribution-valued function $U_{\lambda}$ in (2.5) corresponding to $S=\sin (t|x|) /|x|, t>0$ fixed, is given by

$$
\left.U_{\lambda}(\xi)=\left.\mathcal{D}_{L^{\infty}\left(\mathbb{R}_{x}^{n}\right)}\langle 1(x), \sin (t|x|)| \xi\right|^{\lambda-n+2}|x|^{-\lambda-1} J_{\lambda}(|x| \cdot|\xi|)\right\rangle_{\mathcal{D}_{L^{1}}^{\prime}}\left(\mathbb{R}_{x}^{n}\right) \hat{\otimes} S^{\prime}\left(\mathbb{R}_{\xi}^{n}\right) .
$$

If $\operatorname{Re} \lambda>n-1$, then

$$
\xi \mapsto \sin (t|x|)|\xi|^{\lambda-n+2}|x|^{-\lambda-1} J_{\lambda}(|x| \cdot|\xi|)
$$

is a continuous function with values in $L^{1}\left(\mathbb{R}_{x}^{n}\right)$. Under this assumption on $\lambda$, we therefore obtain that $U_{\lambda} \in \mathcal{C}\left(\mathbb{R}_{\xi}^{n}\right)$ is given by

$$
\begin{equation*}
U_{\lambda}(\xi)=\left|\mathbb{S}^{n-1}\right| \cdot|\xi|^{\lambda-n+2} \int_{0}^{\infty} \rho^{n-\lambda-2} \sin (t \rho) J_{\lambda}(\rho|\xi|) \mathrm{d} \rho \tag{2.8}
\end{equation*}
$$

(b) If $n \geq 3$ is odd, then [29, I, 13.11, p. 67] yields the following for $\operatorname{Re} \lambda>n-1$ :

$$
\begin{aligned}
U_{\lambda}(\xi) & =\left|\mathbb{S}^{n-1}\right| \cdot(-1)^{(n-1) / 2}|\xi|^{\lambda-n+2}\left(\frac{\partial}{\partial t}\right)^{n-2} \int_{0}^{\infty} \rho^{-\lambda} \cos (t \rho) J_{\lambda}(\rho|\xi|) \mathrm{d} \rho \\
& =\left|\mathbb{S}^{n-1}\right| \cdot(-1)^{(n-1) / 2}|\xi|^{\lambda-n+2}\left(\frac{\partial}{\partial t}\right)^{n-2}\left[\frac{\sqrt{\pi}}{(2|\xi|)^{\lambda}} \cdot \chi^{-1 / 2+\lambda}\left(|\xi|^{2}-t^{2}\right)\right] .
\end{aligned}
$$

The distributions $\chi^{-1 / 2+\lambda}\left(|\xi|^{2}-t^{2}\right) \in \mathcal{S}^{\prime}\left(\mathbb{R}_{\xi}^{n}\right)$ depend $\mathcal{C}^{\infty}$ on $t>0$ and hence the last formula holds by analytic continuation for each $\lambda \in \mathbb{C}$ and $t>0$. This implies

$$
\begin{aligned}
E & =(2 \pi)^{-n} Y(t) \mathcal{F}_{x}\left(\frac{\sin (t|x|)}{|x|}\right)=(2 \pi)^{-n} Y(t) \cdot 2^{n / 2-1} \Gamma\left(\frac{n}{2}\right) \cdot U_{n / 2-1}(x) \\
& =\frac{Y(t)}{\Gamma\left(\frac{n-1}{2}\right)(-4 \pi)^{(n-1) / 2}} \partial_{t}^{n-2}\left[Y(|x|-t)|x|^{2-n}\left(|x|^{2}-t^{2}\right)^{(n-3) / 2}\right] \\
& =\frac{1}{(n-2)!\left|\mathbb{S}^{n-1}\right|} \partial_{t}^{n-2}\left[Y(t-|x|)|x|^{2-n}\left(t^{2}-|x|^{2}\right)^{(n-3) / 2}\right]
\end{aligned}
$$

in accordance with [31, Lemma 3.3.5, p. 218 (for $k=1$ )].
(c) For even $n$ and $\operatorname{Re} \lambda>n-1$ we obtain the following from (2.8):

$$
\begin{equation*}
U_{\lambda}(\xi)=\left|\mathbb{S}^{n-1}\right| \cdot(-1)^{n / 2-1}|\xi|^{\lambda-n+2}\left(\frac{\partial}{\partial t}\right)^{n-2} \int_{0}^{\infty} \rho^{-\lambda} \sin (t \rho) J_{\lambda}(\rho|\xi|) \mathrm{d} \rho \tag{2.9}
\end{equation*}
$$

The integral in (2.9) is absolutely convergent for $\operatorname{Re} \lambda>1$ and yields a continuous function of $t$ and $\xi$ depending analytically on $\lambda$. However, this integral is more complicated than the one in the case of odd $n$ (see [15, Eq. 6.699.1, p. 747]), and we proceed differently. By applying once more [35, Prop. 4, p. 41], we obtain from Theorem 2.1 also a formula for the partial Fourier transform, i.e., the distribution-valued function

$$
U: \mathbb{C} \longrightarrow \mathcal{S}^{\prime}\left(\mathbb{R}_{t, \xi}^{m+n}\right): \lambda \longmapsto U_{\lambda}(t, \xi)=\mathcal{S}^{\prime}\left(\mathbb{R}_{t}^{m}\right) \hat{\otimes} \mathcal{S}^{\prime}\left(\mathbb{R}_{x}^{n}\right)\left\langle S(t, x), K_{\lambda}(x, \xi)\right\rangle_{\mathcal{S}\left(\mathbb{R}_{x}^{n}\right) \hat{\otimes} \mathcal{S}^{\prime}\left(\mathbb{R}_{\xi}^{n}\right)}
$$

is entire for $S(t, x) \in \mathcal{S}^{\prime}\left(\mathbb{R}_{t, x}^{m+n}\right)$ and $\left(\mathcal{F}_{x} S\right)(t, \xi)=2^{n / 2-1} \Gamma\left(\frac{n}{2}\right) \cdot U_{n / 2-1}$ if $S$ is radially symmetric with respect to $x$. Hence, putting $m=1, S(t, x)=(2 \pi)^{-n} Y(t) \sin (t|x|) /|x|$ and assuming $n$ even with $n \geq 6$ we can insert $\lambda=\frac{n}{2}-1$ into (2.9), and we obtain by analytic continuation

$$
E=(2 \pi)^{-n / 2} Y(t) \cdot \partial_{t}^{n-2}\left[|\xi|^{-n / 2+1} \int_{0}^{\infty} \rho^{-n / 2+1} J_{-n / 2+1}(\rho|\xi|) \sin (t \rho) \mathrm{d} \rho\right]
$$

(Note that $J_{-k}(s)=(-1)^{k} J_{k}(s)$ for $k \in \mathbb{N}$ and $s \in \mathbb{R}$. Let us also mention that the last formula can be deduced as well for $n=2$ or $n=4$ upon using a further differentiation with respect to $t$.)

In order to evaluate the last integral, let us assume $\xi \neq 0$ and consider the analytic distribution-valued function

$$
T:\left\{v \in \mathbb{C} ; \operatorname{Re} v>-\frac{1}{2}\right\} \longrightarrow \mathcal{S}^{\prime}\left(\mathbb{R}_{\rho}^{1}\right): v \longmapsto T_{v}(\rho)=Y(\rho) \rho^{v} J_{v}(\rho|\xi|) .
$$

By means of the series expansion of the Bessel function, we infer that $T$ can be analytically continued to $\mathbb{C} \backslash\left(-\frac{1}{2}-\mathbb{N}_{0}\right)$ having simple poles in $-\frac{1}{2}-\mathbb{N}_{0}$ (see also [22, 2.4, p. 193]). Furthermore,

$$
T:\left\{v \in \mathbb{C} ; \operatorname{Re} v<-\frac{1}{2}, v \notin-\frac{1}{2}-\mathbb{N}\right\} \longrightarrow \mathcal{D}_{L^{1}}^{\prime}\left(\mathbb{R}_{\rho}^{1}\right)
$$

is also well-defined and analytic. A classical formula (see [29, II, 13.9, p. 164]) furnishes

$$
\begin{equation*}
\int_{0}^{\infty} \rho^{\nu} J_{v}(\rho|\xi|) \sin (t \rho) \mathrm{d} \rho=\sqrt{\pi}(2|\xi|)^{v} \chi^{-1 / 2-v}\left(t^{2}-|\xi|^{2}\right) \tag{2.10}
\end{equation*}
$$

if $-1<\operatorname{Re} v<-\frac{1}{2}$. By analytic continuation we deduce from (2.10)

$$
\begin{aligned}
& \int_{0}^{\infty} \rho^{-n / 2+1} J_{-n / 2+1}(\rho|\xi|) \sin (t \rho) \mathrm{d} \rho \\
& \quad= \mathcal{D}_{L^{\infty}\left(\mathbb{R}_{\rho}^{1}\right)}\left\langle\sin (t \rho), \rho^{-n / 2+1} J_{-n / 2+1}(\rho|\xi|)\right\rangle_{\mathcal{D}_{L^{1}}^{\prime}\left(\mathbb{R}_{\rho}^{1}\right)} \\
& \quad=\sqrt{\pi}(2|\xi|)^{-n / 2+1} \chi^{(n-3) / 2}\left(t^{2}-|\xi|^{2}\right),
\end{aligned}
$$

and this yields

$$
\begin{equation*}
E=\frac{1}{(n-2)!\left|\mathbb{S}^{n-1}\right|} \partial_{t}^{n-2}\left[Y(t-|x|)|x|^{2-n}\left(t^{2}-|x|^{2}\right)^{(n-3) / 2}\right] \tag{2.11}
\end{equation*}
$$

exactly as in the case of odd $n$. Let us mention that a unified deduction of (2.11) independent of the parity of $n$ is given in [31, Lemma 3.3.5, p. 218].

For some examples (and in particular for the application to the EPD-operator in Sect. 3), we need a different family of kernels $\tilde{K}_{\lambda}$. Of course, they are chosen such that $\tilde{K}_{n / 2-1}$ again coincides with $J_{n / 2-1}(|x| \cdot|\xi|) /(|x| \cdot|\xi|)^{n / 2-1}$.

Theorem 2.3 Let $K_{\lambda}$ be defined as in Theorem 2.1 and assume $\lambda \in \mathbb{C} \backslash(-\mathbb{N})$. Then the kernel

$$
\tilde{K}_{\lambda}(x, \xi)=|x|^{2 \lambda-n+2} \cdot K_{\lambda}(x, \xi)=(|x| \cdot|\xi|)^{\lambda-n+2} J_{\lambda}(|x| \cdot|\xi|)
$$

belongs to $\mathcal{O}_{C}^{\prime}\left(\mathbb{R}_{x}^{n}\right) \hat{\otimes} \mathcal{S}^{\prime}\left(\mathbb{R}_{\xi}^{n}\right)$ and it depends therein holomorphically on $\lambda \in \mathbb{C} \backslash(-\mathbb{N})$.
If $S \in \mathcal{S}_{r}^{\prime}\left(\mathbb{R}^{n}\right)$ and $S$ is $\mathcal{C}^{\infty}$ in a neighborhood of 0 , then $S(x) \cdot \tilde{K}_{\lambda}(x, \xi)$ belongs to $\mathcal{D}_{L^{1}}^{\prime}\left(\mathbb{R}_{x}^{n}\right) \hat{\otimes} \mathcal{S}^{\prime}\left(\mathbb{R}_{\xi}^{n}\right)$ and depends therein holomorphically on $\lambda \in \mathbb{C} \backslash(-\mathbb{N})$. Finally,

$$
\mathcal{F} S=2^{n / 2-1} \Gamma\left(\frac{n}{2}\right) \cdot \tilde{U}_{n / 2-1}
$$

where the function

$$
\begin{equation*}
\tilde{U}: \mathbb{C} \backslash(-\mathbb{N}) \rightarrow \mathcal{S}^{\prime}\left(\mathbb{R}_{\xi}^{n}\right): \lambda \mapsto \tilde{U}_{\lambda}(\xi)=\mathcal{D}_{L^{\infty}}\left\langle 1(x), S(x) \cdot \tilde{K}_{\lambda}(x, \xi)\right\rangle_{\mathcal{D}_{L^{1}}^{\prime}\left(\mathbb{R}_{x}^{n}\right) \hat{\otimes} \mathcal{S}^{\prime}\left(\mathbb{R}_{\xi}^{n}\right)} \tag{2.12}
\end{equation*}
$$

is holomorphic.
Proof From Theorem 2.1, we infer then

$$
\tilde{K}_{\lambda}(x, \xi)=|x|^{2 \lambda-n+2} \cdot K_{\lambda}(x, \xi) \in \mathcal{S}^{\prime}\left(\mathbb{R}_{x}^{n}\right) \cdot\left(\mathcal{S}\left(\mathbb{R}_{x}^{n}\right) \hat{\otimes} \mathcal{S}^{\prime}\left(\mathbb{R}_{\xi}^{n}\right)\right) \subset \mathcal{O}_{C}^{\prime}\left(\mathbb{R}_{x}^{n}\right) \hat{\otimes} \mathcal{S}^{\prime}\left(\mathbb{R}_{\xi}^{n}\right)
$$

due to $\mathcal{S} \cdot \mathcal{S}^{\prime} \subset \mathcal{O}_{C}^{\prime}$ and [34, Prop. 20 bis, p. 70]. Furthermore, $\tilde{K}_{\lambda}$ depends holomorphically on $\lambda \in \mathbb{C} \backslash(-\mathbb{N})$ since the same holds for $|x|^{2 \lambda-n+2}$. (One can show that $\tilde{K}_{\lambda}$ has simple poles at $\lambda=-k, k \in \mathbb{N}$. E.g., it holds $\operatorname{Res}_{\lambda=-1} \tilde{K}_{\lambda}=\frac{1}{2}\left|\mathbb{S}^{n-1}\right|^{2} \cdot \delta(x, \xi)$.)

Analogously, also the distribution-valued function

$$
\begin{aligned}
\lambda \mapsto S(x) \cdot \tilde{K}_{\lambda}(x, \xi)=\left(S(x)|x|^{2 \lambda-n+2}\right) \cdot K_{\lambda}(x, \xi) & \in \mathcal{O}_{C}^{\prime}\left(\mathbb{R}_{x}^{n}\right) \hat{\otimes} \mathcal{S}^{\prime}\left(\mathbb{R}_{\xi}^{n}\right) \\
& \subset \mathcal{D}_{L^{1}}^{\prime}\left(\mathbb{R}_{x}^{n}\right) \hat{\otimes} \mathcal{S}^{\prime}\left(\mathbb{R}_{\xi}^{n}\right)
\end{aligned}
$$

is holomorphic on $\mathbb{C} \backslash(-\mathbb{N})$. Finally, $\mathcal{F} S=2^{n / 2-1} \Gamma\left(\frac{n}{2}\right) \cdot \tilde{U}_{n / 2-1}$ since $\tilde{K}_{n / 2-1}=K_{n / 2-1}$. This completes the proof.

Example 2.4 Let us illustrate the difference of the representations for $\mathcal{F} S$ in Theorem 2.1 and Theorem 2.3, respectively, by considering again the forward fundamental solution $E$ of the wave operator $\partial_{t}^{2}-\Delta_{n}$. If $S=\sin (t|x|) /|x|, t>0$ fixed, then (2.12) in Theorem 2.3 yields

$$
\left.\tilde{U}_{\lambda}(\xi)=\left.\mathcal{D}_{L^{\infty}\left(\mathbb{R}_{x}^{n}\right)}\langle 1(x), \sin (t|x|)| \xi\right|^{\lambda-n+2}|x|^{\lambda-n+1} J_{\lambda}(|x| \cdot|\xi|)\right\rangle_{\mathcal{D}_{L^{1}}^{\prime}\left(\mathbb{R}_{x}^{n}\right) \hat{\otimes} S^{\prime}\left(\mathbb{R}_{\xi}^{n}\right)} .
$$

For $-1<\operatorname{Re} \lambda<-\frac{1}{2}$ and $|\xi| \leq N, N \in \mathbb{N}$, the moduli of the functions $f_{\xi}(x)=$ $\sin (t|x|)|\xi|^{-\lambda}|x|^{\lambda-n+1} J_{\lambda}(|x| \cdot|\xi|)$ are bounded, independently of $\xi$, by the integrable function

$$
\begin{aligned}
& g_{\lambda, N}(x)=C_{\lambda}|\sin (t|x|)| \cdot|x|^{2 \operatorname{Re} \lambda-n+1} \cdot(1+N|x|)^{-1 / 2-\operatorname{Re} \lambda} \in L^{1}\left(\mathbb{R}_{x}^{n}\right), \\
& \quad C_{\lambda}=\left\|u^{-\lambda}(1+u)^{1 / 2+\lambda} J_{\lambda}(u)\right\|_{L^{\infty}((0, \infty))} .
\end{aligned}
$$

Therefore, $\tilde{U}_{\lambda}(\xi) \in L_{\mathrm{loc}}^{1}\left(\mathbb{R}_{\xi}^{n}\right)$ for $-1<\operatorname{Re} \lambda<-\frac{1}{2}$, and $\tilde{U}_{\lambda}(\xi)$ is given, for $\xi \neq 0$, by the absolutely convergent integral

$$
\tilde{U}_{\lambda}(\xi)=\left|\mathbb{S}^{n-1}\right| \cdot|\xi|^{\lambda-n+2} \int_{0}^{\infty} \rho^{\lambda} J_{\lambda}(\rho|\xi|) \sin (t \rho) \mathrm{d} \rho
$$

that we have encountered already in (2.10). By analytic continuation, we thus obtain

$$
\tilde{U}_{\lambda}(\xi)=\frac{2^{\lambda+1} \pi^{(n+1) / 2}}{\Gamma\left(\frac{n}{2}\right)}|\xi|^{2 \lambda-n+2} \cdot \chi^{-1 / 2-\lambda}\left(t^{2}-|\xi|^{2}\right), \quad \lambda \in \mathbb{C} \backslash(-\mathbb{N})
$$

Hence we deduce from Theorem 2.3 the following expression for the forward fundamental solution $E$ of $\partial_{t}^{2}-\Delta_{n}$ :

$$
\begin{equation*}
E=(2 \pi)^{-n} Y(t) 2^{n / 2-1} \Gamma\left(\frac{n}{2}\right) \tilde{U}_{n / 2-1}(x)=\frac{1}{2} \pi^{(1-n) / 2} Y(t) \chi^{(1-n) / 2}\left(t^{2}-|x|^{2}\right) \tag{2.13}
\end{equation*}
$$

(As said above, we interpret $E$ as a continuous function of $t$ with values in $\mathcal{S}^{\prime}\left(\mathbb{R}_{x}^{n}\right)$ and vanishing for $t \leq 0$. Furthermore, for $t>0$, the composition $h^{*} T$ of $T=\chi^{(1-n) / 2} \in$ $\mathcal{D}^{\prime}\left(\mathbb{R}^{1}\right)$ with the submersive $\mathcal{C}^{\infty}$ function $h(x)=t^{2}-|x|^{2}$ is well-defined, see (1.4).) The representation of $E$ in (2.13) was given already in [10, Lemme 4.2, p. 317], see also [18, Thm. 6.2.1, (6.2.1)', p. 138] or [31, Ex. 1.6.17, p. 106].

Remark 2.5 Let us eventually observe that we could also employ the kernel

$$
K_{\lambda}^{0}(x, \xi)=\frac{J_{\lambda}(|x| \cdot|\xi|)}{(|x| \cdot|\xi|)^{\lambda}}
$$

for the analytic continuation of the Poisson-Bochner formula, yet only for a restricted class of distributions $S$. In fact, by partial Fourier transformation, it follows that

$$
K_{\lambda}^{0}(x, \xi) \in \mathcal{D}_{L^{\infty}, n}\left(\mathbb{R}_{x}^{n}\right) \hat{\otimes} \mathcal{S}^{\prime}\left(\mathbb{R}_{\xi}^{n}\right) \text { where } \mathcal{D}_{L^{\infty}, n}\left(\mathbb{R}_{x}^{n}\right)=\left(1+|x|^{2}\right)^{-n / 2} \cdot \mathcal{D}_{L^{\infty}}\left(\mathbb{R}_{x}^{n}\right)
$$

Hence, for $S \in \mathcal{D}_{L^{1},-n}^{\prime}\left(\mathbb{R}_{x}^{n}\right)=\left(1+|x|^{2}\right)^{n / 2} \cdot \mathcal{D}_{L^{1}}^{\prime}\left(\mathbb{R}_{x}^{n}\right)$, the function

$$
\begin{equation*}
U^{0}: \mathbb{C} \longrightarrow \mathcal{S}^{\prime}\left(\mathbb{R}_{\xi}^{n}\right): \lambda \longmapsto U_{\lambda}^{0}(\xi)={\mathcal{\mathcal { D } _ { L ^ { 1 } , - n } ^ { \prime }}}^{\left(\mathbb{R}_{x}^{n}\right)}\left\langle S(x), K_{\lambda}^{0}(x, \xi)\right\rangle_{\mathcal{D}_{L^{\infty}, n}\left(\mathbb{R}_{x}^{n}\right) \hat{\otimes} \mathcal{S}^{\prime}\left(\mathbb{R}_{\xi}^{n}\right)} \tag{2.14}
\end{equation*}
$$

is entire and $\mathcal{F} S=2^{n / 2-1} \Gamma\left(\frac{n}{2}\right) \cdot U_{n / 2-1}^{0}$ if $S$ is rotationally invariant.
If $S=\sin (t|x|) /|x|, t>0$ fixed, then the assumption $S \in \mathcal{D}_{L^{1},-n}^{\prime}\left(\mathbb{R}^{n}\right)$ is satisfied and $U_{\lambda}^{0}$ in (2.14) would yield the same representation of $\mathcal{F} S$ as in Example 2.2. If, in contrast, $S=1$, then $S \notin \mathcal{D}_{L^{1},-n}^{\prime}\left(\mathbb{R}^{n}\right)$ and the entire distribution-valued function $U^{0}$ in (2.14) does not exist. Note, however, that $U$ and $\tilde{U}$ in (2.5) and in (2.12), respectively, remain meaningful and yield

$$
U_{\lambda}=\frac{2^{n-\lambda} \pi^{n / 2}|\xi|^{2 \lambda-2 n+2}}{\Gamma\left(\lambda-\frac{n}{2}+1\right)}, \quad \operatorname{Re} \lambda>\frac{n}{2}-1
$$

and

$$
\tilde{U}_{\lambda}=\frac{2^{\lambda+2} \pi^{n} \Gamma(\lambda+1)}{\Gamma\left(\frac{n}{2}\right)^{2}} \delta, \quad \lambda \in \mathbb{C} \backslash(-\mathbb{N})
$$

and $\mathcal{F} 1=2^{n / 2-1} \Gamma\left(\frac{n}{2}\right) U_{n / 2-1}=2^{n / 2-1} \Gamma\left(\frac{n}{2}\right) \tilde{U}_{n / 2-1}=(2 \pi)^{n} \delta$ as expected.

## 3 The fundamental solution of the Euler-Poisson-Darboux operator

Let us turn now to the Euler-Poisson-Darboux operator

$$
\begin{equation*}
P_{\alpha}\left(t, \partial_{t}, \partial_{x}\right)=\partial_{t}^{2}+\frac{2 \alpha+1}{t} \partial_{t}-\Delta_{n} \tag{3.1}
\end{equation*}
$$

acting on the space of distributions defined in the right half-space $\left\{(t, x) \in \mathbb{R}^{n+1} ; t>0\right\}$. Since $P_{\alpha}\left(t, \partial_{t}, \partial_{x}\right)$ is strictly hyperbolic with respect to $t$, it has a unique fundamental solution $E_{\alpha, \tau}(t, x) \in \mathcal{D}^{\prime}\left(\mathbb{R}^{n+1}\right)$ for each $\tau>0$ fulfilling

$$
\begin{equation*}
P_{\alpha}\left(t, \partial_{t}, \partial_{x}\right) E_{\alpha, \tau}(t, x)=\delta_{\tau}(t) \otimes \delta(x) \text { and } \operatorname{supp} E_{\alpha, \tau} \subset\left\{(t, x) \in \mathbb{R}^{n+1} ; t \geq \tau\right\}, \tag{3.2}
\end{equation*}
$$

see [19, Thm. 23.2.2, p. 392], [7, Ch. 6, Thm. and Def. 4.9, p. 379].
Moreover, the strict hyperbolicity of $P_{\alpha}\left(t, \partial_{t}, \partial_{x}\right)$ implies that $E_{\alpha, \tau}$ depends $C^{\infty}$ on $t$ for $t \geq \tau$ and that the support of $E_{\alpha, \tau}$ is contained in the propagation cone $\left\{(t, x) \in \mathbb{R}^{n+1} ; t \geq\right.$ $\tau+|x|\}$. In particular, $E_{\alpha, \tau} \in \mathcal{C}^{\infty}([\tau, \infty)) \hat{\otimes} \mathcal{E}^{\prime}\left(\mathbb{R}^{n}\right)$ and the partial Fourier transform $S_{\alpha, \tau}$ of $E_{\alpha, \tau}$ with respect to $x$ fulfills

$$
S_{\alpha, \tau}=\mathcal{F}_{x}\left(E_{\alpha, \tau}\right) \in \mathcal{C}^{\infty}([\tau, \infty)) \hat{\otimes} \mathcal{O}_{M}\left(\mathbb{R}_{x}^{n}\right)
$$

i.e., $S_{\alpha, \tau}$ is an infinitely differentiable mapping from $[\tau, \infty)$ into $\mathcal{O}_{M}\left(\mathbb{R}^{n}\right)$. By constructing the Green function of the ordinary differential operator $\partial_{t}^{2}+(2 \alpha+1) t^{-1} \partial_{t}+|x|^{2}$, we next derive an explicit representation of $S_{\alpha, \tau}$.

Proposition 3.1 For $\tau>0$ and $\alpha \in \mathbb{C}$, we have

$$
\begin{equation*}
S_{\alpha, \tau}=\frac{\pi}{2} Y(t-\tau) \tau^{\alpha+1} t^{-\alpha}\left[-N_{\alpha}(\tau|x|) J_{\alpha}(t|x|)+J_{\alpha}(\tau|x|) N_{\alpha}(t|x|)\right] . \tag{3.3}
\end{equation*}
$$

[Here $N_{\alpha}, \alpha \in \mathbb{C}$, denote the Bessel functions of the second kind.]
Proof Upon Fourier transform with respect to $x$, (3.2) yields

$$
\begin{equation*}
\left(\partial_{t}^{2}+\frac{2 \alpha+1}{t} \partial_{t}+|x|^{2}\right) S_{\alpha, \tau}(t, x)=\delta_{\tau}(t) . \tag{3.4}
\end{equation*}
$$

This ordinary differential equation arises by substitution from Bessel's equation, and the vector space of its homogeneous solutions is generated by $t^{-\alpha} J_{\alpha}(t|x|)$ and $t^{-\alpha} N_{\alpha}(t|x|)$, see [21, C, 2.162, (9), p. 440] or [15, 8.491.6, p. 971].

Equation (3.4) implies that $S_{\alpha, \tau}$ has the two initial values $S_{\alpha, \tau}(\tau, x)=0$ and $\left(\frac{\mathrm{d}}{\mathrm{d} t} S_{\alpha, \tau}\right)(\tau, x)=1$. If therefore

$$
S_{\alpha, \tau}=Y(t-\tau)\left[C_{1} t^{-\alpha} J_{\alpha}(t|x|)+C_{2} t^{-\alpha} N_{\alpha}(t|x|)\right],
$$

then the constants $C_{1}, C_{2}$ are determined by the following system of linear equations:

$$
0=C_{1} \tau^{-\alpha} J_{\alpha}(\tau|x|)+C_{2} \tau^{-\alpha} N_{\alpha}(\tau|x|),
$$

$$
1=C_{1} \partial_{\tau}\left(\tau^{-\alpha} J_{\alpha}(\tau|x|)\right)+C_{2} \partial_{\tau}\left(\tau^{-\alpha} N_{\alpha}(\tau|x|)\right) .
$$

The Wronskian determinant

$$
W(\tau)=\operatorname{det}\left(\begin{array}{cc}
\tau^{-\alpha} J_{\alpha}(\tau|x|) & \tau^{-\alpha} N_{\alpha}(\tau|x|) \\
\partial_{\tau}\left(\tau^{-\alpha} J_{\alpha}(\tau|x|)\right) & \partial_{\tau}\left(\tau^{-\alpha} N_{\alpha}(\tau|x|)\right)
\end{array}\right)
$$

fulfills $W(\tau)=C \tau^{-2 \alpha-1}$, see [21, A, 17.1, p. 72], and the power series of $J_{\alpha}$ and $N_{\alpha}$ yield $C=\frac{2}{\pi}$. Hence

$$
\binom{C_{1}}{C_{2}}=\frac{\pi}{2} \tau^{2 \alpha+1}\binom{-\tau^{-\alpha} N_{\alpha}(\tau|x|)}{\tau^{-\alpha} J_{\alpha}(\tau|x|)}
$$

and this furnishes formula (3.3).
Formula (3.3) is essentially equivalent with Eq. (5.5) in [8, p. 332], see also [6, pp. 30-32].
In order to evaluate the Fourier transform $E_{\alpha, \tau}=(2 \pi)^{-n} \mathcal{F}_{x}\left(S_{\alpha, \tau}\right)$, we use analytic continuation of the Poisson-Bochner formula according to Theorem 2.3. The representation of the fundamental solution $E_{\alpha, \tau}$ in Proposition 3.2 below coincides with formula (7.4) in [10, Lemme 7.1, p. 327] taking into account [10, (5.9), p. 319]. This result of S. Delache and J. Leray was also obtained (with the same method of proof) in [4, Thm. 2.1, (2.16), p. 501]. Our deduction of $E_{\alpha, \tau}$ is different and seems to be new.

Let us mention that an earlier appearance of this fundamental solution in the form of a "Riemann function" can be found in [40, p. 361, last line]. In fact, $E_{\alpha, \tau}=(\tau / t)^{1 / 2+\alpha}$. $U^{2}(0, \tau ; x, t)$ where, in Young's formula, $m=n+1, \lambda=\alpha, \Gamma=(t-\tau)^{2}-|x-\xi|^{2}$, $H_{m}(2)=2 \pi^{(n-1) / 2} \Gamma\left(\frac{3-n}{2}\right)$ and hence
$U^{2}(0, \tau ; x, t)=\frac{1}{2} \pi^{(1-n) / 2} \chi^{(1-n) / 2}\left((t-\tau)^{2}-|x|^{2}\right){ }_{2} F_{1}\left(\frac{1}{2}-\alpha, \frac{1}{2}+\alpha ; \frac{3-n}{2} ;-\frac{(t-\tau)^{2}-|x|^{2}}{4 \tau t}\right)$.
(Note that $H_{n}(\alpha+2)$ is defined erroneously in [40, p. 357] and that the equation $\lambda=(1-k) / 2$ in [40, p. 361] should read $\lambda=(k-1) / 2$.)

Proposition 3.2 Let $\tau>0$ and assume that $n=1$ or $n \in \mathbb{N}$ is even. Then the fundamental solution $E_{\alpha, \tau} \in \mathcal{S}^{\prime}\left(\mathbb{R}_{t, x}^{n+1}\right)$ of the Euler-Poisson-Darboux operator $P_{\alpha}\left(t, \partial_{t}, \partial_{x}\right)=\partial_{t}^{2}+$ $\frac{2 \alpha+1}{t} \partial_{t}-\Delta_{n}$, i.e., the unique solution of (3.2), is given by

$$
\begin{align*}
E_{\alpha, \tau}(t, x)= & \frac{1}{2} \pi^{(1-n) / 2}\left(\frac{\tau}{t}\right)^{1 / 2+\alpha} Y(t-\tau) \chi^{(1-n) / 2}\left((t-\tau)^{2}-|x|^{2}\right) \\
& \times{ }_{2} F_{1}\left(\frac{1}{2}-\alpha, \frac{1}{2}+\alpha ; \frac{3-n}{2} ;-\frac{(t-\tau)^{2}-|x|^{2}}{4 \tau t}\right) \tag{3.5}
\end{align*}
$$

[As said above, we interpret $E_{\alpha, \tau}$ as a continuous function of $t$ with values in $\mathcal{E}^{\prime}\left(\mathbb{R}_{x}^{n}\right)$ and vanishing for $t \leq \tau$. Furthermore, for $t>\tau$, the composition $h^{*} T$ of $T=\chi^{(1-n) / 2} \in \mathcal{D}^{\prime}\left(\mathbb{R}^{1}\right)$ with the submersive $\mathcal{C}^{\infty}$ function $h(x)=(t-\tau)^{2}-|x|^{2}$ is well-defined, see (1.4), and so is the multiplication with the $\mathcal{C}^{\infty}$ function given by ${ }_{2} F_{1}$.]

Proof (a) In order to apply Theorem 2.3, let us first check that $S_{\alpha, \tau}(t, x)$ in (3.5) is a $\mathcal{C}^{\infty}$ function of $x$ in a neighborhood of 0 . In fact, for $\alpha \in \mathbb{C} \backslash \mathbb{Z}$, we have

$$
\begin{aligned}
S_{\alpha, \tau}(t, x) & =\frac{\pi}{2} Y(t-\tau) \tau^{\alpha+1} t^{-\alpha}\left[-N_{\alpha}(\tau|x|) J_{\alpha}(t|x|)+J_{\alpha}(\tau|x|) N_{\alpha}(t|x|)\right] \\
& =\frac{\pi \tau^{\alpha+1} Y(t-\tau)}{2 \sin (\alpha \pi) t^{\alpha}}\left[J_{-\alpha}(\tau|x|) J_{\alpha}(t|x|)-J_{\alpha}(\tau|x|) J_{-\alpha}(t|x|)\right]
\end{aligned}
$$

which is an analytic function of $|x|^{2}$, see [15, Equs. 8.402 and 8.403.1, p. 951]. The same holds for $\alpha \in \mathbb{Z}$ by taking limits. Therefore, Theorem 2.3 applies and yields (for $t>\tau$ fixed) $\mathcal{F}_{x} S_{\alpha, \tau}=2^{n / 2-1} \Gamma\left(\frac{n}{2}\right) \tilde{U}_{n / 2-1}$, where

$$
\begin{equation*}
\tilde{U}_{\lambda}(\xi)=\mathcal{D}_{L^{\infty}\left(\mathbb{R}_{x}^{n}\right)}\left\langle 1(x), S_{\alpha, \tau}(t, x) \cdot(|x| \cdot|\xi|)^{\lambda-n+2} J_{\lambda}(|x| \cdot|\xi|)\right\rangle_{\mathcal{D}_{L^{1}}^{\prime}\left(\mathbb{R}_{x}^{n}\right) \hat{\otimes} \mathcal{S}^{\prime}\left(\mathbb{R}_{\xi}^{n}\right)} \tag{3.6}
\end{equation*}
$$

for $\lambda \in \mathbb{C} \backslash(-\mathbb{N})$.
(b) If $-1<\operatorname{Re} \lambda<-\frac{1}{2}$, then we see, analogously as in Example 2.4, which is the special case of $\alpha=-\frac{1}{2}$, that $\tilde{U}_{\lambda}(\xi) \in L_{\mathrm{loc}}^{1}\left(\mathbb{R}^{n}\right)$ and that the evaluation in (3.6) furnishes an absolutely convergent integral for $\xi \neq 0$. Hence we obtain the following for $-1<\operatorname{Re} \lambda<-\frac{1}{2}$, $0<\tau<t$ and $\xi \neq 0$ fixed:

$$
\begin{aligned}
\left(\tilde{U}_{\lambda}(\xi)=\right. & \frac{\pi^{n / 2+1} \tau^{\alpha+1}}{\Gamma\left(\frac{n}{2}\right) t^{\alpha}}|\xi|^{\lambda-n+2} \\
& \times \int_{0}^{\infty} \rho^{\lambda+1}\left[-N_{\alpha}(\tau \rho) J_{\alpha}(t \rho)+J_{\alpha}(\tau \rho) N_{\alpha}(t \rho)\right] \cdot J_{\lambda}(\rho|\xi|) \mathrm{d} \rho .
\end{aligned}
$$

According to [28, 10.51, 10.52, p. 93] (see also [1, (7.1), p. 45]), this integral yields

$$
\tilde{U}_{\lambda}(\xi)=\frac{\sqrt{2} \pi^{(n+1) / 2} \tau^{\alpha-\lambda}}{\Gamma\left(\frac{n}{2}\right) t^{\alpha+\lambda+1}} Y(t-\tau-|\xi|)|\xi|^{2 \lambda-n+2}\left(u^{2}-1\right)^{-\lambda / 2-1 / 4} p_{-1 / 2+\alpha}^{1 / 2+\lambda}(u)
$$

where

$$
u=\frac{t^{2}+\tau^{2}-|\xi|^{2}}{2 t \tau} \geq 1 \text { for } t-\tau \geq|\xi| .
$$

The Legendre function $p_{-1 / 2+\alpha}^{1 / 2+\lambda}$ can be expressed by Gauß' hypergeometric function, see [28, p. 279], and this leads to

$$
\begin{align*}
\tilde{U}_{\lambda}(x)= & \frac{2^{\lambda+1} \pi^{(n+1) / 2}}{\Gamma\left(\frac{n}{2}\right)}\left(\frac{\tau}{t}\right)^{1 / 2+\alpha}|x|^{2 \lambda-n+2} \chi^{-1 / 2-\lambda}\left((t-\tau)^{2}-|x|^{2}\right) \\
& \times{ }_{2} F_{1}\left(\frac{1}{2}-\alpha, \frac{1}{2}+\alpha ; \frac{1}{2}-\lambda ;-\frac{(t-\tau)^{2}-|x|^{2}}{4 \tau t}\right) . \tag{3.7}
\end{align*}
$$

(c) Since $\tilde{U}_{\lambda}$ is holomorphic in $\mathbb{C} \backslash(-\mathbb{N})$, formula (3.7) persists by analytic continuation for those $\lambda \in \mathbb{C} \backslash(-\mathbb{N})$ where the right-hand side is, in this form, well-defined, i.e., for $\lambda \notin \frac{1}{2}+\mathbb{N}_{0}$. Note that ${ }_{2} F_{1}(a, b ; c ; u)$ has poles for $c \in-\mathbb{N}_{0}$. In particular, if $n=1$ or $n$ is even and $\lambda=$ $\frac{n}{2}-1$, then the hypergeometric series is regular at $c=\frac{1}{2}-\lambda=\frac{3-n}{2}$ and we obtain straightforwardly formula (3.5) in the proposition by using $E_{\alpha, \tau}=(2 \pi)^{-n} 2^{n / 2-1} \Gamma\left(\frac{n}{2}\right) \tilde{U}_{n / 2-1}(x)$. This completes the proof.

For odd space dimensions $n \geq 3$, Eq. (3.7) in the proof of Proposition 3.2 is still valid, but the analytic continuation into the point $\lambda=\frac{n}{2}-1$ needs more care due to the simple poles of the function $c \mapsto{ }_{2} F_{1}(a, b ; c ; u)$ at $c=-k \in \mathbb{N}_{0}$. The representation of the fundamental solution $E_{\alpha, \tau}$ in Proposition 3.3 below coincides with formula (7.4) in [10, Lemme 7.1, p. 327] taking into account [10, (5.10), p. 320]. This result of S. Delache and J. Leray was also obtained (with the same method of proof) in [4, Thm. 2.2, (2.18), p. 502].

Proposition 3.3 Let $\tau>0$ and assume that $n \geq 3$ and $n \in \mathbb{N}$ is odd. Then the fundamental solution $E_{\alpha, \tau} \in \mathcal{S}^{\prime}\left(\mathbb{R}_{t, x}^{n+1}\right)$ of the Euler-Poisson-Darboux operator $P_{\alpha}\left(t, \partial_{t}, \partial_{x}\right)=\partial_{t}^{2}+$
$\frac{2 \alpha+1}{t} \partial_{t}-\Delta_{n}$, i.e., the unique solution of (3.2), is given by

$$
\begin{align*}
E_{\alpha, \tau}= & \frac{1}{2 \pi^{k+1}}\left(\frac{\tau}{t}\right)^{1 / 2+\alpha} Y(t-\tau)\left[\sum_{j=0}^{k} c_{j} \delta^{(k-j)}\left((t-\tau)^{2}-|x|^{2}\right)\right. \\
& \left.+c_{k+1} Y\left((t-\tau)^{2}-|x|^{2}\right)_{2} F_{1}\left(\frac{n}{2}-\alpha, \frac{n}{2}+\alpha ; \frac{n+1}{2} ; u\right)\right] \tag{3.8}
\end{align*}
$$

where

$$
k=\frac{n-3}{2}, \quad c_{j}=\frac{\left(\frac{1}{2}-\alpha\right)_{j}\left(\frac{1}{2}+\alpha\right)_{j}}{j!(-4 \tau t)^{j}}, \quad u=-\frac{(t-\tau)^{2}-|x|^{2}}{4 \tau t}
$$

and the Pochhammer symbol $(w)_{j}$ is as in the introduction.
[As in Proposition 3.2, $E_{\alpha, \tau}$ is interpreted as a continuous function of $t$ with values in $\mathcal{E}^{\prime}\left(\mathbb{R}_{x}^{n}\right)$ and vanishing for $t \leq \tau$.]

Proof (a) From the series expansion in [15, 9.100, p. 1039] and the transformation formula [15, 9.131.1, p. 1043], we see that

$$
f: \mathbb{C} \backslash\left(-\mathbb{N}_{0}\right) \longrightarrow \mathcal{E}((-\infty, 1)): c \longmapsto f_{c}(u)={ }_{2} F_{1}(a, b ; c ; u)
$$

is a holomorphic function for fixed $a, b \in \mathbb{C}$ having at most simple poles in $c=-k, k \in \mathbb{N}_{0}$. The formula [15, 9.101.1, p. 1039] furnishes the residues:

$$
\begin{align*}
\operatorname{Res}_{c=-k} f_{c}(u) & =\lim _{c \rightarrow-k}(c+k)_{2} F_{1}(a, b ; c ; u) \\
& =\left(\operatorname{Res}_{c=-k} \Gamma(c)\right) \cdot \lim _{c \rightarrow-k} \frac{{ }_{2} F_{1}(a, b ; c ; u)}{\Gamma(c)} \\
& =\frac{(-1)^{k}(a)_{k+1}(b)_{k+1} u^{k+1}}{k!(k+1)!}{ }_{2} F_{1}(a+k+1, b+k+1 ; k+2 ; u) . \tag{3.9}
\end{align*}
$$

Here we have used Pochhammer's symbol $(w)_{j}$ as defined in Sect. 1 .
Furthermore, the Taylor series of $\mathrm{Pf}_{c=-k} f_{c}(z)$ up to order $k$ is given by

$$
\begin{equation*}
\operatorname{Pf}_{c=-k} f_{c}(u)=\sum_{j=0}^{k} \frac{(a)_{j}(b)_{j}}{(-k)_{j} j!} u^{j}+O\left(u^{k+1}\right) . \tag{3.10}
\end{equation*}
$$

(b) Let us next investigate the holomorphic distribution-valued function
$T:\left\{c \in \mathbb{C} ; 2 c-3 \notin \mathbb{N}_{0}\right\} \longrightarrow \mathcal{S}^{\prime}\left(\mathbb{R}_{x}^{n}\right): c \longmapsto T_{c}=2^{-c}|x|^{-2 c-n+3} \chi^{c-1}\left((t-\tau)^{2}-|x|^{2}\right)$ for fixed $0<\tau<t$ and near a point $c=-k, k \in \mathbb{N}_{0}$. On the one hand, $\chi^{-k-1}=\delta^{(k)}$ (see [18, (3.2.17)', p. 74]) implies

$$
\begin{equation*}
T_{-k}=2^{k}|x|^{2 k-n+3} \delta^{(k)}\left((t-\tau)^{2}-|x|^{2}\right) . \tag{3.11}
\end{equation*}
$$

On the other hand, for $\operatorname{Re} c>0$, we have

$$
\frac{\mathrm{d} \chi^{c-1}(s)}{\mathrm{d} c}=\frac{\mathrm{d}}{\mathrm{~d} c}\left(\frac{Y(s) s^{c-1}}{\Gamma(c)}\right)=\frac{Y(s) s^{c-1}}{\Gamma(c)} \cdot[\log s-\psi(c)]
$$

and in particular

$$
\left.\frac{\mathrm{d} \chi^{c-1}(s)}{\mathrm{d} c}\right|_{c=1}=Y(s) \cdot[\log s-\psi(1)] \in L_{\mathrm{loc}}^{1}\left(\mathbb{R}_{s}^{1}\right)
$$

For $k \in \mathbb{N}_{0}$, this implies

$$
\begin{aligned}
\left.\frac{\mathrm{d} \chi^{c-1}(s)}{\mathrm{d} c}\right|_{c=-k} & =\partial_{s}^{k+1}\left(\left.\frac{\mathrm{~d} \chi^{c-1}(s)}{\mathrm{d} c}\right|_{c=1}\right)=\partial_{s}^{k+1}(Y(s) \cdot[\log s-\psi(1)]) \\
& =\partial_{s}^{k}\left(s_{+}^{-1}-\psi(1) \delta\right)=(-1)^{k} k!s_{+}^{-k-1}-\psi(k) \delta^{(k)},
\end{aligned}
$$

see $\left[31\right.$, p. 50] for the last equation. Here $s_{+}^{-k-1}=\operatorname{Pf}_{c=-k-1} Y(s) s^{c} \in \mathcal{S}^{\prime}\left(\mathbb{R}_{s}^{1}\right)$ and $\psi(k)=$ $\Gamma^{\prime}(k) / \Gamma(k)$. So eventually, we obtain

$$
\begin{align*}
\left.\frac{\mathrm{d} T_{c}}{\mathrm{~d} c}\right|_{c=-k}= & 2^{k}|x|^{2 k-n+3}\left[(-1)^{k} k!s_{+}^{-k-1}\left((t-\tau)^{2}-|x|^{2}\right)\right. \\
& \left.-\left(\psi(k)+\log \left(2|x|^{2}\right)\right) \cdot \delta^{(k)}\left((t-\tau)^{2}-|x|^{2}\right)\right] . \tag{3.12}
\end{align*}
$$

(c) As we have observed above, the distribution-valued function $\tilde{U}_{\lambda}$ is holomorphic in $\mathbb{C} \backslash(-\mathbb{N})$ and thus is regular at $\lambda=\frac{n}{2}-1$. In fact, setting $n=2 k+3, k \in \mathbb{N}_{0}, a=\frac{1}{2}-\alpha$, $b=\frac{1}{2}+\alpha, c=\frac{1}{2}-\lambda$ and $u=-\left[(t-\tau)^{2}-|x|^{2}\right] /(4 \tau t)$ and assuming $0<\tau<t$ fixed, we obtain

$$
\operatorname{Res}_{\lambda=n / 2-1} \tilde{U}_{\lambda}=-\frac{2^{3 / 2} \pi^{(n+1) / 2}}{\Gamma\left(\frac{n}{2}\right)}\left(\frac{\tau}{t}\right)^{1 / 2+\alpha} T_{-k} \cdot \operatorname{Res}_{c=-k} f_{c}(u)=0
$$

since $T_{-k}=2^{k} \delta^{(k)}(-4 \tau t u)$ by (3.11) and $\operatorname{Res}_{c=-k} f_{c}(u)$ vanishes of order $k+1$ at $u=0$ by (3.9).

Similarly, using [30, Prop. 1.6.3, p. 28] we conclude that

$$
\begin{aligned}
\tilde{U}_{n / 2-1}= & \operatorname{Pf}_{\lambda=n / 2-1} \tilde{U}_{\lambda}=\frac{2^{3 / 2} \pi^{(n+1) / 2}}{\Gamma\left(\frac{n}{2}\right)}\left(\frac{\tau}{t}\right)^{1 / 2+\alpha} \operatorname{Pf}_{c=-k}\left[T_{c}(x) \cdot f_{c}(u)\right] \\
= & \frac{2^{3 / 2} \pi^{(n+1) / 2}}{\Gamma\left(\frac{n}{2}\right)}\left(\frac{\tau}{t}\right)^{1 / 2+\alpha}\left[\left.\frac{\mathrm{d} T_{c}}{\mathrm{~d} c}\right|_{c=-k}\right. \\
& \left.\cdot \operatorname{Res}_{c=-k} f_{c}(u)+T_{-k}(x) \cdot \operatorname{Pf}_{c=-k} f_{c}(u)\right] .
\end{aligned}
$$

Equations (3.9/3.12) yield

$$
\begin{aligned}
& \left.\frac{\mathrm{d} T_{c}}{\mathrm{~d} c}\right|_{c=-k} \cdot \operatorname{Res}_{c=-k} f_{c}(u)=\frac{(-1)^{k+1}\left(\frac{1}{2}-\alpha\right)_{k+1}\left(\frac{1}{2}+\alpha\right)_{k+1}}{2^{k+2}(k+1)!(\tau t)^{k+1}} \\
& \quad \times Y\left((t-\tau)^{2}-|x|^{2}\right){ }_{2} F_{1}\left(\frac{n}{2}-\alpha, \frac{n}{2}+\alpha ; \frac{n+1}{2} ; u\right)
\end{aligned}
$$

and Eqs. (3.10/3.11) yield

$$
T_{-k}(x) \cdot \operatorname{Pf}_{c=-k} f_{c}(u)=2^{k} \delta^{(k)}\left((t-\tau)^{2}-|x|^{2}\right) \cdot \sum_{j=0}^{k} \frac{(a)_{j}(b)_{j}}{(-k)_{j} j!} u^{j} .
$$

Upon using the identity

$$
s^{j} \delta^{(k)}(s)=(-k)_{j} \delta^{(k-j)}(s) \in \mathcal{D}^{\prime}\left(\mathbb{R}_{s}^{1}\right), \quad j, k \in \mathbb{N}_{0}, j \leq k,
$$

this results in

$$
T_{-k}(x) \cdot \operatorname{Pf}_{c=-k} f_{c}(u)=2^{k} \cdot \sum_{j=0}^{k} \frac{(a)_{j}(b)_{j}}{j!(-4 \tau t)^{j}} \delta^{(k-j)}\left((t-\tau)^{2}-|x|^{2}\right)
$$

Finally, we make use of $E_{\alpha, \tau}=(2 \pi)^{-n} 2^{n / 2-1} \Gamma\left(\frac{n}{2}\right) \tilde{U}_{n / 2-1}$ in order to conclude the representation of $E_{\alpha, \tau}$ in (3.8). This completes the proof.

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