# Analytic Dependence is an Unnecessary Requirement in Renormalization of Locally Covariant QFT 

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#### Abstract

Finite renormalization freedom in locally covariant quantum field theories on curved spacetime is known to be tightly constrained, under certain standard hypotheses, to the same terms as in flat spacetime up to finitely many curvature dependent terms. These hypotheses include, in particular, locality, covariance, scaling, microlocal regularity and continuous and analytic dependence on the metric and coupling parameters. The analytic dependence hypothesis is somewhat unnatural, because it requires that locally covariant observables (which are simultaneously defined on all spacetimes) depend continuously on an arbitrary metric, with the dependence strengthened to analytic on analytic metrics. Moreover the fact that analytic metrics are globally rigid makes the implementation of this requirement at the level of local $*$-algebras of observables rather technically cumbersome. We show that the conditions of locality, covariance, scaling and a naturally strengthened microlocal spectral condition, are actually sufficient to constrain the allowed finite renormalizations equally strongly, thus eliminating both the continuity and the somewhat unnatural analyticity hypotheses. The key step in the proof uses the Peetre-Slovák theorem on the characterization of (in general non-linear) differential operators by their locality and regularity properties.


## 1 Introduction

Perturbative ultraviolet renormalization of locally covariant quantum field theories in (globally hyperbolic) curved spacetime is a well established topic of algebraic quantum field theory, especially for scalar fields [5, 6, 13, 14]. It essentially deals with two classes of objects: Wick polynomials and time ordered Wick polynomials. Exactly as in flat spacetime, these objects can be considered as the building blocks of the whole renormalization procedure. Smeared versions of Wick polynomials, of their time ordered products and of their derivatives

[^0]generate an algebra $\mathcal{W}(M, \mathbf{g})$, for a given spacetime $(M, \mathbf{g})$, enlarged in a controlled way from the algebra of products of smeared linear fields. This enlarged algebra then includes physically fundamental observables, such as the stressenergy tensor, which is necessary, for instance, to evaluate the energy densities and fluxes of physical processes in curved spacetimes like particle creation or Hawking radiation. The stress-energy tensor is also needed to compute the back reaction of the quantum matter on the background geometry.

This paper deals only with Wick polynomials, or more precisely just Wick powers with all results easily extended to all Wick polynomials by linearity, though the presented results could in principle be adapted to deal also with their derivatives and their time ordered products. In curved spacetime, Wick polynomials have to satisfy stronger locality and covariance requirements than in flat spacetime. These requirements are conveniently stated in the language of category theory introduced in [7], which we also use here. We should stress, though, that the categorical language primarily serves to compress somewhat long lists of hypotheses into concise statements. Existence of locally covariant Wick polynomials and their time ordered products was established in the seminal works of Hollands and Wald, respectively in [13] and [14]. It is well known that, in flat spacetime, time ordered Wick polynomials are not uniquely defined. This fact survives the passage to curved spacetime. However, unlike in flat spacetime, the absence of a preferred reference state means that Wick polynomials are themselves not uniquely defined. The ambiguities involved with the definition of these two classes of fields are physically interpreted as finite renormalizations or renormalization counterterms, upon adopting the natural locally covariant generalization of Epstein-Glaser approach to renormalization.

Exactly as in flat spacetime, each fixed type of (either Wick or time-ordered) polynomial admits a finite-dimensional class of independent counterterms. In curved spacetime, this class is much larger than in Minkowski space, because of the possible dependence of counterterms on background curvature. While this class may no longer be finite-dimensional, it is still finitely generated or quasi-finite-dimensional in a precise sense, because the counterterms may depend only polynomially on the curvature scalars up to a certain dimension. This remarkable result, in the case of Wick polynomials, presented in [13, Thm. 5.1] and summarized before the statement of our Theorem 3.1, is arrived at by imposing severe constraints on Wick polynomials in addition to those of locality and covariance. These requirements are of various kinds. Some arise from heuristic properties of quantum free fields, e.g., Hermiticity and commutation relations. Other requirements concern microlocal features which, loosely speaking, extend to curved spacetime the structure of Fourier transforms of the relevant Green functions on Minkowski space. Another requirement regards the behaviour of Wick polynomials under a rescaling of the metric and the parameters $m^{2}$ and $\xi$ of the free theory, which describe the field's mass ${ }^{1}$ and its coupling to the curvature. Finally there are the technically delicate requirements of continuous and analytic dependence on the metric. The two latter requirements play a crucial role in [13] in their proof of the strong restrictions on possible finite renormalization counterterms that was mentioned above.

The main difficulty with defining a suitable notion of the continuous depen-

[^1]dence of an element of the algebra $\mathcal{W}(M, \mathbf{g})$ on the metric $\mathbf{g}$ (and the other parameters $m^{2}$ and $\xi$ ) is that, continuously changing the metric $\mathbf{g} \mapsto \mathbf{g}^{\prime}$, the whole algebra $\mathcal{W}(M, \mathbf{g})$ changes correspondingly and algebras $\mathcal{W}(M, \mathbf{g})$ and $\mathcal{W}\left(M, \mathbf{g}^{\prime}\right)$ associated with different metrics are not canonically isomorphic. Therefore even just stating the condition of continuous dependence on $\mathbf{g}$ requires some finesse. Locality can be turned into an advantage in this context [13]. One may restrict attention to metric variations in a spacetime region $O \subset M$ with compact closure. If $\mathbf{g}$ agrees with $\mathbf{g}^{\prime}$ outside $O$, essentially exploiting a suitable version of the time slice axiom, it is possible to naturally identify an element of $\mathcal{W}(M, \mathbf{g})$ with a corresponding element in $\mathcal{W}\left(M, \mathbf{g}^{\prime}\right)$, when both are supported in $O$. Hence, a local version of the continuity requirement can be imposed by means of this canonical identification.

The requirement of analytic dependence is even trickier to state. It is argued in [13] that analytic dependence is necessary because the remaining requirements would not be able to rule out the undesirable infinite family of non-polynomial in curvature counterterms that were considered in [25]. There is an important subtle technical issue that arises in stating this analytic dependence condition. The way followed for stating the continuity dependence requirement in a local region $O$ faces here an insurmountable obstruction: analytic metrics are rigid and if they coincide outside $O$ they must coincide also in $O$. The ingenious but cumbersome strategy elaborated in [13] makes use of a special class of Hadamard states over the considered algebras. Since no local analytic variations of the metric are possible, they consider a joint analytic family $\mathbf{g}^{(s)}$ of the metric on $O$ and a corresponding analytic family of quasifree Hadamard states $\omega^{(s)}$ on $\mathcal{W}\left(M, \mathbf{g}^{(s)}\right)$. Then they require that the distributions obtained by composing $\omega^{(s)}$ with the local Wick polynomials (or their time ordered products) varies analytically with $s$ in a suitable analytic and microlocal sense (see the discussion starting on p. 311 in [13]).

Continuous and analytic dependence on the parameters $m^{2}$ and $\xi$ is there treated similarly, with both parameters taken to be functions on $M$, rather than just constants, at intermediate stages of the arguments.

The main result of this work establishes that the technically cumbersome and somewhat unnatural analytic dependence requirement is in fact not necessary to achieve the classification theorem [13, Thm. 5.1]. Our classification result, Theorem 3.1, is essentially identical, though it is slightly more general because it allows smooth (rather than just analytic) dependence on the dimensionless curvature coupling $\xi$. In our proof, we make no use of the continuous and analytic dependence requirements of [13]. Instead, leaving all the other requirements on Wick powers the same, we appeal to a strengthened (more precisely, parametrized) version of the microlocal spectrum condition, which we believe is natural from both the physical and geometrical points of view. This modification of the axioms on Wick powers is sufficient to achieve the desired classification. Also, echoing the original (rather implicit) arguments of [13], we believe that it is very likely that our version of the axioms are satisfied by the standard locally covariant Hadamard parametrix prescription for explicitly constructing Wick powers, which would mean that our classification results are non-vacuous. However, we leave a detailed verification of this claim to future work.

The key tool exploited in our proof is a theorem that characterizes (generally non-linear) differential operators in terms of their locality and regularity
properties. This theorem, known as the Peetre-Slovák theorem (or sometimes the non-linear Peetre theorem), in its most elementary version (Proposition 2.2; see also Appendix A for a more general statement) states the following: any map $D$, that associates smooth sections $\psi: M \rightarrow E$ of a bundle $E \rightarrow M$ to smooth sections $D[\psi]: M \rightarrow F$ of another bundle $F \rightarrow M$ in such a way that $D[\psi](x)$ depends only on the germ of $\psi$ at $x$ for any point $x \in M$, is necessarily a differential operator of locally bounded order, smoothly depending on its arguments and their derivatives. The $C_{k}$ coefficients that characterize renormalization counterterms of Wick polynomials precisely map sections of the bundle of metrics and parameters, $m^{2}$ and $\xi$, to scalar valued distributions on a spacetime $M$. The microlocal conditions ensure that these distributions are actually smooth functions, while the locality requirement implies that the $C_{k}$ satisfy the hypotheses of Peetre-Slovák's theorem and hence must be differential operators. A combination of the scaling and covariance requirements then shows that the differential order of the $C_{k}$ is globally bounded and that their dependence on the metric, $m^{2}$ and the derivatives of all the parameters is polynomial, with coefficients smoothly depending on $\xi$. Further, covariance also dictates that the derivatives of the metric necessarily group into curvature scalars.

Notably, the analytic dependence requirement is not exploited in establishing the above result. Within the context of our proof, counter terms like $m^{k} F\left(R / m^{2}\right)$, where $R$ is the Ricci scalar and $F$ is any smooth function with strong decay near 0 and $\pm \infty$, as considered in [25], are excluded because they violate the microlocal requirement: there exists a choice of a spacetime $(M, \mathbf{g})$ and of a scalar field $m^{2}$ such that the counterterm is not smooth and hence has non-empty wavefront set and the Wick polynomials modified by adding these counterterms do not satisfy the microlocal requirement (neither the original, nor our strengthened version).

This paper is organized as follows. Our main theorem and its proof are presented in Sect. 3. The proof is somewhat lengthy, but straight forward. It relies on some preliminary definitions and results discussed in Sect. 2. In particular our basic version of Peetre-Slovák's theorem is stated in Sect. 2.3 after a quick summary of elementary facts about jet bundles in Sect. 2.2, where we also introduce some useful coordinate systems. Sect. 2.4 is devoted to introducing our notion of scaling which is more precise but substantially equivalent to the one employed in [13]. However, we are careful to identify two different kinds of scalings (physical and coordinate), which were mixed in [13] by the introduction of Riemann normal coordinates. The remainder of Sect. 2 deals with notions and results, especially on $G L(n)$ representation theory, which are useful for imposing the covariance requirement. After recalling the definition and properties of Wick polynomials, and the more general notion of locally covariant quantum field, with the appropriate categorical language, we state and prove our main result in several steps in Sect. 3. Sect. 4 concludes the paper with a discussion of the results and directions for future work. Appendix A illustrates a more general version of Peetre-Slovák's theorem, which applies to differential operators with parameters.

## 2 Geometry of scaling and general covariance

In this section we discuss some aspects of the geometry of the higher derivatives (jets) of metric and scalar fields under the action of scaling and diffeomorphism transformations. These properties will be crucial in the characterization of finite renormalizations in locally covariant quantum field theory in Sect. 3.

### 2.1 Coordinates on jets

In differential geometry, jets $[18,16]$ are a geometric way of collecting information about higher derivatives of functions (or bundle sections) on manifolds, similar to what the tangent and cotangent bundles do for first derivatives. Jets have an invariant geometric meaning even on manifolds without a preferred metric or connection. Further, a choice of a coordinate chart on a manifold induces a choice of adapted coordinates on the corresponding jet bundle. One advantage of working with jets is that certain calculations are very conveniently performed in such an adapted local coordinate chart, yet also lead to global and geometrically invariant conclusions. Below, we briefly discuss some variations on adapted local coordinate systems on the space of jets of bundle of metrics with some scalar fields.

Consider a smooth map $f: \mathbb{R}^{m} \rightarrow \mathbb{R}^{n}$, such that $f(0)=0$. The germ of $f$ at $0 \in \mathbb{R}^{m}$ is the equivalence class of smooth maps $f^{\prime}: \mathbb{R}^{m} \rightarrow \mathbb{R}^{n}$ that agree with $f$ on some neighborhood of $0 \in \mathbb{R}^{m}$. The $r$-jet of $f$ at $0 \in \mathbb{R}^{m}$ is the equivalence class of all smooth maps $f^{\prime}: \mathbb{R}^{m} \rightarrow \mathbb{R}^{n}$ that have the same Taylor expansion at 0 as $f$ to order $r$, denoted $j_{0}^{r} f$. Obviously, the germ contains more information than a jet of any order. These definitions are clearly local, both on the domain and the target of a smooth map, and are invariant under $C^{\infty}$-changes of coordinates. Thus, these definitions easily translate to maps between smooth finite-dimensional (smooth) manifolds $M, N$ replacing $0 \in \mathbb{R}^{m}$ and $0 \in \mathbb{R}^{n}$, respectively, by generic points $x \in M, y \in N$. In particular, with the said $M$ and $N$, we denote by $J^{r}(M, N)$ the set of all distinct jets $j_{x}^{r} f$ of all smooth maps $f: M \rightarrow N$ for all $x \in M$. Also, if $E \rightarrow N$ is a smooth bundle over $N$, then we denote by $J^{r} E$ or, for emphasis, by $J^{r}(E \rightarrow N) \subset J^{r}(N, E)$ the subset of jets of smooth sections $f: N \rightarrow E$. Both $J^{r}(M, N)$ and $J^{r}(E \rightarrow N)$ can be given structures of smooth manifolds. A fiber $\left(J^{r} E\right)_{x}$ at $x \in N$ is diffeomorphic to $E_{x} \times \mathbb{R}^{s_{r}}$, where $E_{x}$ is the fiber of $E$ and $s_{r}$ counts the components of all (symmetrized) partial derivatives up to order $r$. In fact, by projection onto the target of each jet, $J^{r} E \rightarrow E \rightarrow N$ is an iterated smooth bundle. Given a section $\psi: N \rightarrow E$, we can collect the $r$-jets of $\psi$ over each point of $N$ into a section $j^{r} \psi: N \rightarrow J^{r} E$ called the $r$-jet extension of $\psi$.

Let $\left(x^{a}, v^{i}\right)$ be a local adapted coordinate chart on a bundle $F \rightarrow M$, where ( $x^{a}$ ) serve as coordinates on a domain $U \subseteq M$ and $\left(x^{a}, v^{i}\right)$ serve as trivializing coordinates on the fibers of the domain $V \subseteq F$ over $U$. For example, if $T_{q}^{p} M \rightarrow$ $M$ is the bundle of $(p, q)$-tensors we can choose coordinates $\left(x^{a}, t_{b_{1} \cdots b_{q}}^{a_{1} \cdots a_{p}}\right)$ on the projection pre-image $V$ of $U$, such that a section $\tau: T_{q}^{p} M \rightarrow M$ could locally be written as

$$
\begin{equation*}
\tau(x)=t_{b_{1} \cdots b_{q}}^{a_{1} \cdots a_{p}}(\tau(x)) \mathrm{d} x^{b_{1}} \cdots \mathrm{~d} x^{b_{q}} \frac{\partial}{\partial x^{a_{1}}} \cdots \frac{\partial}{\partial x^{a_{p}}} \tag{1}
\end{equation*}
$$

The local chart $\left(x^{a}, v^{i}\right)$ then induces an adapted coordinate system $\left(x^{a}, v_{A}^{i}\right)$ on
the domain $V^{r} \subseteq J^{r} E$ that is the projection pre-image of $V$ and is diffeomorphic to $V^{r} \cong V \times \mathbb{R}^{s_{r}}$, with $s_{r}$ as discussed above. Each $A=a_{1} \cdots a_{l}$, standing in for an unordered (equivalently, fully symmetrized) collection of base manifold coordinate indices, is a multi-index of size $|A|=l$, with the range $l=0,1, \ldots, r$. The defining property of these coordinates is the identity

$$
\begin{equation*}
v_{A}^{i}\left(j^{r} \psi(x)\right)=\partial_{A} v^{i}(\psi(x))=\frac{\partial}{\partial x^{a_{1}}} \cdots \frac{\partial}{\partial x^{a_{l}}} v^{i}(\psi(x)), \tag{2}
\end{equation*}
$$

for any section $\psi: M \rightarrow F$. Given such a coordinate system, for brevity, we use the notation $\partial_{a}=\partial / \partial x^{a}$ and $\partial_{i}^{A}=\partial / \partial v_{A}^{i}$ for corresponding coordinate vector fields.

### 2.2 Coordinates on jets of metric and scalar fields

If $M$ is a $n$-dimensional smooth manifold, let us now fix the bundle $B M \rightarrow M$ given by the bundle product of the bundle $\dot{S}^{2} T_{*} M$ of (smooth) Lorentzian metric (0,2)-tensors over $M$ and the trivial bundle $\mathbb{R} \times M \rightarrow M$ of (smooth) scalar fields over $M$. Let us denote the sections of this bundle by $(\mathbf{g}, \xi): M \rightarrow B M$. There are several local coordinate systems on $J^{r} B M$, of various merits, which we discuss below.

Covariant coordinates. Given a local coordinate chart ( $x^{a}$ ) on $U \subseteq M$, we define the corresponding adapted coordinates $\left(x^{a}, g_{a b}, z\right)$ on $V \subseteq B M$, which in turn induce the covariant coordinates

$$
\begin{equation*}
\left(x^{a}, g_{a b, A}, z_{A}\right) \quad \text { on } V^{r} \subseteq J^{r} B M \tag{3}
\end{equation*}
$$

Notice that only $n(n+1) / 2$ components of $g_{a b}$ take part in the above coordinates, in view of the symmetry of the metric.

Contravariant coordinates. Recall that a Lorentzian metric g: $M \rightarrow \dot{S}^{2} T^{*} M$ is invertible and hence defines a section $\mathbf{g}^{-1}: M \rightarrow \dot{S}^{2} T M$. The components of the inverse metric can be extracted by functions $g^{a b}$ defined on all of $V \subseteq$ $B M$, such that $g^{a b}\left(\mathbf{g}^{-1}(x)\right)=g_{a b}(\mathbf{g}(x))$, which induce the functions $g_{A}^{a b}$ on $V^{r}$ that satisfy $g_{A}^{a b}\left(j^{r} \mathbf{g}(x)\right)=\partial_{A} g^{a b}(\mathbf{g}(x))$. Then, using the notation $g^{A B}=$ $g^{a_{1} b_{1}} \cdots g^{a_{l} b_{l}}$, for $|A|=|B|=l$, we define the following functions

$$
\begin{equation*}
g=\left|\operatorname{det} g_{a b}\right|, \quad g^{a b, A}=g^{A B} g_{B}^{a b}, \quad z^{A}=g^{A B} z_{A} \tag{4}
\end{equation*}
$$

where, by invertibility of Lorentzian metrics, the function $g^{-1}$ is well defined on all of $V^{r}$, since $g=\left|\operatorname{det} g_{a b}\right|$ is never zero. These functions make up the alternative set of local contravariant coordinates

$$
\begin{equation*}
\left(x^{a}, g^{a b, A}, z^{A}\right) \quad \text { on } V^{r} \subseteq J^{r} B M \tag{5}
\end{equation*}
$$

with the caveat that as the set of functions $\left(g, g^{a b}\right)$ is only functionally independent up to the identity $g^{-1}=\left|\operatorname{det} g^{a b}\right|$, for instance, one of the contravariant coordinates $g^{a b}$ can be replaced by $g$. These coordinates have convenient scaling properties that will be exploited in Sect. 2.4.

Rescaled contravariant coordinates. Another coordinate system that we introduce on $V^{r} \subseteq J^{r} B M$, the rescaled contravariant coordinates, is a suitable rescaling of the previous one. Namely, we introduce various factors of $g^{\alpha}$ in the latter coordinates ( $n$ being the dimension of $M$ ):

$$
\begin{equation*}
\left(x^{a}, g, g^{-\frac{1}{n}} g_{a b}, g^{\frac{1}{n}+\frac{1}{n}|A|} g^{a b, A}, g^{\frac{s}{2 n}+\frac{1}{n}|A|} z^{A}\right) \tag{6}
\end{equation*}
$$

where one of the $n(n+1) / 2$ functions $g^{-\frac{1}{n}} g_{a b}$ is omitted and replaced by $g$. This is because the functions $g^{-\frac{1}{n}} g_{a b}$ are not functionally independent because of the relation $\left|\operatorname{det} g^{-\frac{1}{n}} g^{a b}\right|=1$.

Curvature coordinates. Recall also that, given a Lorentzian metric g, we can always define the corresponding covariant derivative, or Levi-Civita connection, $\nabla$ and the Riemann tensor $\mathbf{R}$. Using well known formulas, we can define functions $\Gamma_{b c}^{a}$ and $\bar{R}_{a b c d}$ on $V^{r} \subseteq J^{r} B M$ that correspond to the coordinate components of the Christoffel symbols and the fully covariant Riemann tensor. Define also the fully contravariant tensor $\mathbf{S}$ with components

$$
\begin{equation*}
\bar{S}^{a b c d}=g^{a a^{\prime}} g^{b b^{\prime}} \bar{R}_{a^{\prime}}\left(c_{b^{\prime}}^{d)}=g^{a b, c d}-g^{b(c, d) a}-g^{a(d, c) b}+g^{c d, a b}+\right.\text { l.o.t } \tag{7}
\end{equation*}
$$

where l.o.t stands for terms that involve only coordinates of lower derivative order. Finally, let $\Gamma_{b c, A}^{a}$ denote the components of the coordinate $\partial_{A}$ derivatives of $\Gamma_{b c}^{a}$, let $\bar{S}^{a b c d, A}$ denote the components of the symmetrized contravariant $\nabla^{A}=\nabla^{\left(a_{1}\right.} \cdots \nabla^{\left.a_{l}\right)}$ derivatives of $\mathbf{S}$, and let $\bar{z}^{A}$ the components of the symmetrized contravariant $\nabla^{A}$ derivatives of the scalar field $\xi$. It is wellknown $[15,2]^{2}$ that

$$
\begin{equation*}
\left(x^{a}, g_{a b}, \Gamma_{(b c, A)}^{a}, \bar{S}^{a b(c d, A)}, \bar{z}^{A}\right) \tag{8}
\end{equation*}
$$

also defines a coordinate system on $V^{r} \subseteq J^{r} B M$, which we shall call curvature coordinates. Note that the barred coordinate functions correspond to components of fully contravariant tensors. These coordinate have convenient transformation properties under diffeomorphisms that will be exploited in Sect. 2.5.

Rescaled curvature coordinates. The final coordinate system that we introduce on $V^{r} \subseteq J^{r} B M$, the rescaled curvature coordinates, merges some of the properties of the systems $\left(x^{a}, g^{-\frac{1}{n}} g_{a b}, g^{\frac{1}{n}+\frac{1}{n}|A|} g^{a b, A}, g^{\frac{s}{2 n}+\frac{1}{n}|A|} z^{A}\right)$ and $\left(x^{a}, g_{a b}, \Gamma_{(b c, A)}^{a}, \bar{S}^{a b(c d, A)}, \bar{z}^{A}\right)$. Namely, we again introduce various factors of $g$ in the curvature coordinates:

$$
\begin{equation*}
\left(x^{a}, g, g^{-\frac{1}{n}} g_{a b}, \Gamma_{(b c, A)}^{a}, g^{\frac{3}{n}+\frac{1}{n}|A|} \bar{S}^{a b(c d, A)}, g^{\frac{s}{2 n}+\frac{1}{n}|A|} \bar{z}^{A}\right) \tag{9}
\end{equation*}
$$

where, again, one of the $n(n+1) / 2$ functions $g^{-\frac{1}{n}} g_{a b}$ is omitted and replaced by $g$.

### 2.3 Locality and the Peetre-Slovák theorem

It is well known that linear differential operators have the property that they are support non-increasing. The powerful, original result of Peetre [19, 20] shows that this property is sufficient to characterize them in the context of $C^{\infty}$ differential geometry. A similar characterization holds even for non-linear differential operators [23, 16, 17], a version of which we present below.

Before proceeding, we need a robust geometric notion of what a differential operator is. Often, differential operators are defined by their expressions in coordinate charts. Any such definition is necessarily coordinate dependent and must be checked to agree on chart overlaps. On the other hand, we can give a

[^2]coordinate independent and global definition of differential operators using jets and the $r$-jet extension map $j^{r}$ defined earlier in Sect. 2.2.

Given a smooth bundle $E \rightarrow N$, recall that the $r$-jet extension acts as a map $j^{r}: \Gamma(E \rightarrow N) \rightarrow \Gamma\left(J^{r} E \rightarrow N\right)$, where as usual $\Gamma(G \rightarrow L)$ denotes the space of smooth sections of the bundle $G \rightarrow L$. For our purposes, the map $j^{r}$ will serve as a universal differential operator of order $r$ in the following sense.

Definition 2.1. Let $E \rightarrow N$ and $F \rightarrow M$ be smooth bundles, and consider a $\operatorname{map} D: \Gamma(E) \rightarrow \Gamma(F)$.
(a) $D$ is a differential operator of globally bounded order if there exists an integer $r \geq 0$, the order, and a smooth function $d: J^{r}(E \rightarrow N) \rightarrow$ $F$, considered as a bundle map (i.e., fiber preserving), such that for any section $\psi \in \Gamma(E)$ we have an associated section of the form $D[\psi]=d \circ$ $j^{r} \psi \in \Gamma(F)$.
(b) $D$ is a differential operator of locally bounded order if it satisfies a similar condition locally. Namely, for any point of $y \in N$ and section $\phi \in \Gamma(E)$, there exists a neighborhood $U \subseteq N$ of $y$ with compact closure, together with an integer $r \geq 0$, an open neighborhood $V^{r} \subseteq J^{r}(E \rightarrow N)$ of $j^{r} \phi(U)$ projecting onto $U$, and a smooth function $d: V^{r} \rightarrow F$ that respects the projections $V^{r} \rightarrow U$ and $F \rightarrow M$, such that $D[\psi](x)=d \circ j^{r} \psi(x)$ for any $x \in U$ and any $\psi \in \Gamma(E)$ with $j^{r} \psi(U) \subset V^{r}$.

If $E \rightarrow M$ and $F \rightarrow M$ are vector bundles over the same base manifold $M$ and $D: \Gamma(E) \rightarrow \Gamma(F)$ is a linear map such that $\phi(x)=D[\psi](x)$ depends only on the germ of $\psi$ at $x \in M$ then it is clear that $D$ will be support nonincreasing. Elementary reasoning shows that a linear, support non-increasing map will also only depend on germs. So, another way to rephrase the Peetre theorem for linear differential operators is as follows, where the dependence on the germ replaces the support non-increasing property.

Proposition 2.1 (Linear Peetre's Theorem [19, 20]). Let $E \rightarrow M$ and $F \rightarrow M$ be vector bundles and $D: \Gamma(E) \rightarrow \Gamma(F)$ a linear map such that $\phi(x)=D[\psi](x)$ depends only on the germ of $\psi$ at $x \in M$. Then $D$ is a linear differential operator of locally bounded order (with smooth coefficients in view of the above definition).

In other words, despite the fact that germs potentially contain much more information that jets, such a linear map that depends only on germs in fact sees only jets.

Phrased as above, in terms of germs, the hypotheses of Peetre's theorem are immediately adaptable to the case when the map $D$ is non-linear and acts on sections of (non-vector) smooth bundles. We will only require an additional regularity ${ }^{3}$ hypothesis.

Definition 2.2. Given smooth bundles $E \rightarrow N$ and $F \rightarrow M$, a smooth ( $k$ dimensional) family of sections of $E \rightarrow N$ is a smooth section of the pullback bundle $\pi^{*} E \rightarrow \mathbb{R}^{k} \times N$ (cf. Eq. (71)), where $\pi: \mathbb{R}^{k} \times N \rightarrow N$ is the projection

[^3]onto the second factor, and similarly for families of sections $F \rightarrow M$. A map $D: \Gamma(E \rightarrow N) \rightarrow \Gamma(F \rightarrow M)$ is regular if it maps smooth families of sections to smooth families of sections. A smooth family $\sigma: \mathbb{R}^{k} \times N \rightarrow \pi^{*} E$ is called a compactly supported variation if there exists a compact subset $O \subset N$ such that $\sigma$ is constant along the $\mathbb{R}^{k}$ factor on the complement $\mathbb{R}^{k} \times N \backslash \pi^{-1}(O)$. The map $D$ is weakly regular if it maps smooth compactly supported variations to smooth compactly supported variations.

Proposition 2.2 (Peetre-Slovák's Theorem). Let $E \rightarrow M$ and $F \rightarrow M$ be smooth bundles and $D: \Gamma(E) \rightarrow \Gamma(F)$ a map such that $\phi(x)=D[\psi](x)$ depends only on the germ of $\psi$ at $x \in M$. If in addition $D$ is weakly regular, then it is a (non-linear) differential operator of locally bounded order.

This proposition will be sufficient for our purposes. However, in the standard literature [23], [16, § 19], this result is stated in much greater generality. In fact, that level of generality can obscure the meaning and significance of the theorem. Though, it should be noted that a simplified statement of the theorem, essentially identical to the one above, together with a straight-forward self-contained proof recently appeared in [17]. Note that these standard references usually require regularity instead of weak regularity, but a slight modification of the proof given in [17] makes it clear that only weak regularity is necessary. This point is discussed in Appendix A. Also in Appendix A, we briefly introduce the language needed to state a more general version, Proposition A.1. The above simpler version becomes a special case of Proposition A. 1 once it is trivially checked that $D$ is id-local, where id: $M \cong M$ is the identity map. The more general result given in Appendix A serves two purposes. The first is that it introduces the language in which the Peetre-Slovák theorem and its proof appear in the standard literature $[16, \S 19]$, which also refers to it as the non-linear Peetre theorem. Second, it allows the treatment of differential operators with parameters. For instance, later in Sect. 3, we treat the mass $m^{2}$ of a scalar field and its coupling to curvature $\xi$ as space-time dependent background fields. If they were treated as necessarily spacetime-constant parameters, we would need to substitute Proposition A. 1 for the simpler Proposition 2.2 in the proof of our main Theorem 3.1.

### 2.4 Physical scaling

Referring to the already introduced bundle $B M \rightarrow M$, sections $(\mathbf{g}, \xi) \in \Gamma(B M)$ consist of a smooth Lorentzian metric $\mathbf{g}$ and a smooth scalar field $\xi$ on $M$. We consider the following scaling transformation $(\mathbf{g}, \xi) \mapsto\left(\lambda^{-2} \mathbf{g}, \lambda^{s} \xi\right)$ on sections. We call this transformation a physical scaling, in contrast to a different kind of scaling to be introduced in Sect. 2.5. We will need the following rather general recursive definition, where $\mathbb{R}^{+}:=(0,+\infty)$,

Definition 2.3. Consider a linear representation of the multiplicative group $\mathbb{R}^{+}$on a vector space $W$, written as $W \ni F \mapsto F_{\lambda} \in W$, for every $\lambda \in \mathbb{R}^{+}$.
(a) An element $F \in W$ is said to have homogeneous degree $k \in \mathbb{R}$ if

$$
\begin{equation*}
F_{\lambda}=\lambda^{k} F \quad \text { for all } \lambda \in \mathbb{R}^{+} . \tag{10}
\end{equation*}
$$

(b) An element $F \in W$ is said to have almost homogeneous degree $k \in \mathbb{R}$ and order $l \in \mathbb{N}$ if $l \geq 0$ is an integer such that (the sum over $j$ is omitted if $l=0$ )

$$
\begin{equation*}
F_{\lambda}=\lambda^{k} F+\lambda^{k} \sum_{j=1}^{l}\left(\log ^{j} \lambda\right) G_{j}, \quad \text { for all } \lambda \in \mathbb{R}^{+} \tag{11}
\end{equation*}
$$

and for some $G_{j} \in W$ depending on $F$, which have respectively almost homogeneous degree $k$ and order $l-j$.

The definition is recursive, with higher orders defined in terms of lower ones. Clearly, an element that is almost homogeneous of order $l=0$ is simply homogeneous.
Remark 2.1. Besides almost homogeneous, other common names found in the literature include poly-homogeneous, associated homogeneous and even quasi associated homogeneous. We are mostly interested in the case when $W$ is some function space and the action of $\mathbb{R}^{+}$is induced from an action on the domain of the functions. Reference [22] reviews several definitions leading to this class of functions and lists relevant earlier works. In the context of distribution theory, the terminology of associated homogeneous is prevalent and goes back to the seminal references [10, § 1.4] and [11, Ch.I §4]. Our Definition 2.3 coincides with [22, Def. 5.2].

The physical scaling transformation on the sections $\Gamma(B M)$ can be implemented by post-composing a section with a bundle map $B M \rightarrow B M$ :

$$
\begin{equation*}
B M \ni(p, \mathbf{g}(p), z(p)) \mapsto\left(p, \lambda^{-2} \mathbf{g}(p), \lambda^{s} z(p)\right) \in B M \tag{12}
\end{equation*}
$$

where the real $\lambda \in \mathbb{R}^{+}$defines the scaling transformation. This representation of the multiplicative group $\mathbb{R}^{+}$is globally defined, however this global action can be written in adapted local coordinates, as discussed in Sect. 2.2, and looks like

$$
\begin{equation*}
x^{a} \mapsto x^{a}, \quad g_{a b} \mapsto \lambda^{-2} g_{a b}, \quad z \mapsto \lambda^{s} z . \tag{13}
\end{equation*}
$$

This global transformation lifts to a global transformation of the jet bundle $J^{r} B M$. In the corresponding induced local coordinates, the lifted action reads

$$
\begin{equation*}
g_{a b, A} \mapsto \lambda^{-2} g_{a b, A}, \quad z_{A} \mapsto \lambda^{s} z_{A} \tag{14}
\end{equation*}
$$

We are interested in applying Definition 2.3 to $W=C^{\infty}\left(J^{r} B M\right)$ and the $\mathbb{R}^{+}$ action induced by the lift of physical scalings to $J^{r} B M$. Moreover, we will need to consider also smaller domains $V^{r} \subseteq J^{r} B M$ for these functions, with $V^{r}$ themselves not invariant under physical scalings. Thus, it is more convenient to refer to the infinitesimal version of these transformations, which are effected by the following vector field

$$
\begin{equation*}
e=-2 g_{a b, A} \partial^{a b, A}+s z_{A} \partial_{z}^{A} \tag{15}
\end{equation*}
$$

in the sense that the induced action on scalar functions on $J^{r} B M$ satisfies

$$
\begin{equation*}
\left.\frac{d}{d \lambda}\right|_{\lambda=1} F_{\lambda}=\mathcal{L}_{e} F \tag{16}
\end{equation*}
$$

(In the rest of the paper if $X$ is a vector field on $J^{r} B M, \mathcal{L}_{X}$ denotes the standard Lie derivative so that, in particular $\mathcal{L}_{X}(F):=X(F)$ if $F: V^{r} \subseteq J^{r} B M \rightarrow \mathbb{R}$
is a smooth function.) Notice that, as the physical scaling transformation is globally defined, $e$ turns out to be globally defined on $J^{r} B M$ and (15) is just its expression in local coordinates. We have a first elementary result stated within the following lemma. We will essentially show later that the converse implication holds as well.

Lemma 2.3. A smooth function $F: J^{r} B M \rightarrow \mathbb{R}$ that has almost homogeneous degree $k$ and order $l$, according to Definition 2.3, when the action $F \rightarrow F_{\lambda}$ is the one induced by physical scaling transformations, satisfies the following local infinitesimal version

$$
\begin{equation*}
\left(\mathcal{L}_{e}-k\right)^{l+1} F=0 \tag{17}
\end{equation*}
$$

Proof. It is sufficient to make use of equation (16) and recall the obvious identity $(\lambda d / d \lambda-k)^{l+1} \lambda^{k} \log ^{l} \lambda=0$.

This lemma is essentially a restatement of Theorem 5.2 and Remark 5.1 from [22]. It now permits us to give a definition of almost homogeneity under infinitesimal scaling for functions defined on subsets of jets of dimensionful bundles. The advantage of using infinitesimal scaling is that the domain on which it is defined need not actually be invariant under finite scaling.

Definition 2.4. A smooth function $F: V^{r} \subseteq J^{r} B M \rightarrow \mathbb{R}$, where $V^{r}$ is an open subset which may coincide with all of $J^{r} B M$, is said to have almost homogeneous degree $k \in \mathbb{R}$ and order $l \in \mathbb{N}$ (with $l \geq 0$ ) under physical scalings if it satisfies the identity

$$
\begin{equation*}
\left(\mathcal{L}_{e}-k\right)^{l+1} F=0 . \tag{18}
\end{equation*}
$$

If $l=0, F$ is said to have homogeneous degree $k \in \mathbb{R}$.
To investigate the local structure of $F$ above we initially use an open subset $V^{r}$ equipped with the contravariant coordinates $\left(x^{a}, g^{a b, A}, z^{A}\right)$ introduced in Sect. 2.2. In these coordinates, finite and infinitesimal physical scalings take the form

$$
\begin{gather*}
x^{a} \mapsto x^{a}, \quad g \mapsto \lambda^{-2 n} g, \quad g^{a b, A} \mapsto \lambda^{2+2|A|} g^{a b, A}, \quad z^{A} \mapsto \lambda^{s+2|A|} z^{A},  \tag{19}\\
e=(2+2|A|) g^{a b, A} \partial_{a b, A}+(s+2|A|) z^{A} \partial_{A}^{z}, \tag{20}
\end{gather*}
$$

where we have also described the action of rescaling on $g$ which, as already remarked, can be used as an alternative coordinate in place of one of the $g^{a b}$. As $e$ does not vanish anywhere, $J^{r} B M$ and hence the domain $V^{r}$ are foliated by integral curves of the vector field $e$. Moreover, the identity $\mathcal{L}_{e} g^{-\frac{1}{2 n}}=g^{-\frac{1}{2 n}}$ means that $g$ restricts to a global coordinate on each orbit of $e$. Thus, the level sets of $g$ constitute another foliation of $J^{r} B M$ and $V^{r}$, transverse to the integral curves of $e$. These observations suggest to study the structure of (almost) homogeneous functions of degree $k$ in the rescaled contravariant coordinates

$$
\begin{equation*}
\left(x^{a}, g, g^{-\frac{1}{n}} g_{a b}, g^{\frac{1}{n}+\frac{1}{n}|A|} g^{a b, A}, g^{\frac{s}{2 n}+\frac{1}{n}|A|} z^{A}\right) \tag{21}
\end{equation*}
$$

that were introduced in Sect. 2.2. Note that each of these functions but $g$ is invariant under physical scalings. We have the following result.

Lemma 2.4. Suppose that $V^{r} \subseteq J^{r} B M$ is an open set equipped with either coordinates $\left(x^{a}, g^{a b, A}, z^{A}\right)$ or some other coordinate system introduced in Sect. 2.2, and $F: V^{r} \rightarrow \mathbb{R}$ is a smooth function that has almost homogeneous degree $k$ and order $l$ with respect to physical scalings, as in Definition 2.4. Then there exist homogeneous of degree 0 functions $H_{j}: V^{r} \rightarrow \mathbb{R}$, for $j=0,1, \ldots, l$, such that

$$
\begin{equation*}
F=g^{-\frac{k}{2 n}} \sum_{j=0}^{l} \log ^{j}\left(g^{-\frac{1}{2 n}}\right) H_{j} \tag{22}
\end{equation*}
$$

In particular, using rescaled contravariant coordinates, each $H_{j}$ can be taken independent of $g$ and written in the form

$$
\begin{equation*}
H_{j}=H_{j}\left(x^{a}, g^{-\frac{1}{n}} g_{a b}, g^{\frac{1}{n}+\frac{1}{n}|A|} g^{a b, A}, g^{\frac{s}{2 n}+\frac{1}{n}|A|} z^{A}\right) \tag{23}
\end{equation*}
$$

Proof. In the simplest $l=0$ case, we can define $H=g^{\frac{k}{2 n}} F$ and show that $\mathcal{L}_{e} H=0$ because $\mathcal{L}_{e} g^{-\frac{1}{2 n}}=g^{-\frac{1}{2 n}}$. This means that, in rescaled contravariant coordinates, $H$ is independent of $g$ and hence (23) holds, with $H$ in place of $H_{j}$. Next, the general $l \geq 1$ case can be treated as follows. Let $G:=g^{\frac{k}{2 n}} F$, which implies that $\mathcal{L}_{e}^{l+1} G=g^{\frac{k}{2 n}}\left(\mathcal{L}_{e}-k\right)^{l+1} F=0$. Now, note the identity $\mathcal{L}_{e}^{j} \log ^{j}\left(g^{-\frac{1}{2 n}}\right)=j!$. So, if $H_{l}:=\frac{1}{l!} \mathcal{L}_{e}^{l} G$ and $G_{l-1}:=G-\log ^{l}\left(g^{-\frac{1}{2 n}}\right) H_{l}$, then $\mathcal{L}_{e} H_{l}=0$ and $\mathcal{L}_{e}^{l} G_{l-1}=0$. In other words, starting with $G_{l}=G$, we can recursively define $H_{j}:=\frac{1}{j!} \log ^{j}\left(g^{-\frac{1}{2 n}}\right) \mathcal{L}_{e}^{j} G_{j}$ and $G_{j-1}:=G_{j}-\log ^{j}\left(g^{-\frac{1}{2 n}}\right) H_{j}$, finding $\mathcal{L}_{e} H_{j}=0$ at each step. The procedure stops for $j=0$ when it gives $G_{0}=H_{0}$, so that $G_{j<0}=H_{j<0}=0$, proving (22).

We will also need the following basic result regarding products of vectors with almost homogeneous degree as in Definition 2.3. Due to the generality of Definition 2.3 we must clarify the meaning of product. If $W$ and $W^{\prime}$ are two vector spaces, by a product between them, we mean any fixed bilinear map $W \times W^{\prime} \rightarrow V$, where $V$ is another vector space. If $F \in W$ and $F^{\prime} \in W^{\prime}$ the corresponding element in $V$, their product, will be simply denoted by $F F^{\prime} \in V$.

Lemma 2.5. Referring to Definition 2.3, consider a pair of vector spaces $W, W^{\prime}$ endowed with corresponding representations of $\mathbb{R}^{+}$. Concerning (b) below, assume also that there is a product $W \times W^{\prime} \rightarrow V$ such that (i) $V$ admits a representation of $\mathbb{R}^{+}$and (ii) the map $W \times W^{\prime} \rightarrow V$ is equivariant: $F_{\lambda} F_{\lambda}^{\prime}=\left(F F^{\prime}\right)_{\lambda}$ for $F \in W, F^{\prime} \in W^{\prime}$ and $\lambda \in \mathbb{R}^{+}$.
(a) A linear combination of two elements $F, F^{\prime} \in W$ of almost homogeneous degree $k$ and order $l$ is of almost homogeneous degree $k$ and order $l$.
(b) A product of an element $F \in W$, of almost homogeneous degree $k$ and order $l$, and an element $F^{\prime} \in W^{\prime}$, of almost homogeneous degree $k^{\prime}$ and order $l^{\prime}$, has almost homogeneous degree $k+k^{\prime}$ and order $l+l^{\prime}$.

Proof. Part (a) is trivial, because the defining identity 11 is linear.
We will prove part (b) by double induction on the pair of orders $\left(l, l^{\prime}\right)$.

Consider the identity

$$
\begin{align*}
&\left(F F^{\prime}\right)_{\lambda}=F_{\lambda} F_{\lambda}^{\prime}=\lambda^{k+k^{\prime}} F F^{\prime} \\
&+\lambda^{k+k^{\prime}} \sum_{j=1}^{l}\left(\log ^{j} \lambda\right) G_{j} F^{\prime} \\
&+\lambda^{k+k^{\prime}} \sum_{j^{\prime}=1}^{l^{\prime}}\left(\log ^{j^{\prime}} \lambda\right) F G_{j^{\prime}}^{\prime}  \tag{24}\\
&+\lambda^{k+k^{\prime}} \sum_{j=1}^{l} \sum_{j^{\prime}=1}^{l^{\prime}}\left(\log ^{j+j^{\prime}} \lambda\right) G_{j} G_{j^{\prime}}^{\prime} .
\end{align*}
$$

From this formula, it is clear that, to show that $F F^{\prime}$ has almost homogeneous degree $k+k^{\prime}$ and order $l+l^{\prime}$, it is sufficient to establish that the coefficients of the logarithmic terms, $G_{j} F^{\prime}, F G_{j^{\prime}}^{\prime}$ and $G_{j} G_{j^{\prime}}^{\prime}$, either do not appear or are themselves almost homogeneous of the right degree and order. Thus, to establish the case $\left(l, l^{\prime}\right)$, it is sufficient to have all of the $\left(j, l^{\prime}\right),\left(l, j^{\prime}\right)$ and $\left(j, j^{\prime}\right)$ cases, with $j<l$ and $j^{\prime}<l^{\prime}$, already established. We shall refer to this last remark as the primary inductive step.

The case $\left(l, l^{\prime}\right)=(0,0)$ follows immediately from Eq. (24), since no logarithmic terms appear. Next, we establish the following secondary inductive step. Assuming that, given some $m \geq 0$, all cases ( $l, l^{\prime}$ ) with $l, l^{\prime} \leq m$ hold, then actually all cases $\left(l, l^{\prime}\right)$ with $l, l^{\prime} \leq m+1$ hold as well. To see that, note that the case $(m+1,0)$ holds, because in (24) we need only consider the terms $G_{j} F^{\prime}$, which correspond to the inductively covered cases $(m+1-j, 0)$ with $j \geq 1$. Then, using the primary inductive step, all the cases $\left(m+1, l^{\prime}\right)$ with $1 \leq l^{\prime} \leq m$ follow as well. The cases $(l, m+1)$ with $0 \leq l \leq m$, are completely analogous. Finally, one more appeal to the primary inductive step establishes the case $(m+1, m+1)$.

Iterating the secondary inductive step completes the proof of part (b).

### 2.5 Diffeomorphisms and coordinate scalings

Because the sections $(\mathbf{g}, \xi) \in \Gamma(B M)$ are tensor fields, there is a well defined action of the group $\operatorname{Diff}(M)$ of diffeomorphisms $\chi: M \rightarrow M$ on them by pullback $(\mathbf{g}, \xi) \mapsto\left(\chi^{*} \mathbf{g}, \chi^{*} \xi\right)$. This action of course can be implemented at the level of the bundle itself, $\chi^{*}: B M \rightarrow B M$ and of course lifted to the jet bundle $j^{r} \chi^{*}: J^{r} B M \rightarrow J^{r} B M$. We are interested in the structure of functions $F: J^{r} B M \rightarrow \mathbb{R}$ that are invariant under the action of $\operatorname{Diff}(M)$. We could also consider invariance only under the subgroup $\operatorname{Diff}^{+}(M)$ of orientation preserving diffeomorphisms in an essentially analogous way. For this purpose, it is convenient to make use of the local adapted curvature coordinates $\left(x^{a}, g_{a b}, \Gamma_{(b c, A)}^{a}, \bar{S}^{a b(c d, A)}, \bar{z}^{A}\right)$ on a domain $V^{r} \subseteq J^{r} B M$ defined in Sect. 2.2.

The domain $V^{r}$ itself may not be invariant under $\operatorname{Diff}(M)$, because our coordinates are adapted to a single coordinate chart $\left(x^{a}\right)$ on $U \subseteq M$. On the other hand, having already chosen our coordinate system, we can phrase the requirement that $F: V^{r} \rightarrow \mathbb{R}$ is the restriction of a $\operatorname{Diff}(M)$-invariant function (necessarily defined on a possibly larger $\operatorname{Diff}(M)$-invariant domain) to $V^{r}$ in the following way: (a) $\frac{\partial}{\partial x^{a}} F=0$, where the vector fields $\frac{\partial}{\partial x^{a}}$ are the infinitesimal generators of diffeomorphisms that restrict to coordinate translations on $U$, and (b) the restriction $F_{x}: V_{x}^{r} \subseteq J_{x}^{r} B M \rightarrow \mathbb{R}$ of $F$ to the fiber of $J^{r} B M$ over any one point $x \in M$ is invariant under the action of the subgroup $\operatorname{Diff}(M, x) \subset$
$\operatorname{Diff}(M)$ that fixes $x$. Clearly we can take $V_{x}^{r}$ to be invariant under $\operatorname{Diff}(M, x)$. An immediate simplification based on requirement (a) is that our function is expressible as $F=F_{x}\left(g_{a b}, \Gamma_{(b c, A)}^{a}, \bar{S}^{a b(c d, A)}, \bar{z}^{A}\right)$, that is, it is independent of the base coordinates $\left(x^{a}\right)$. Next, we examine the consequences of requirement (b).

The action of $\operatorname{Diff}(M, x)$ on $r$-jets is not faithful. In fact, it has a large kernel, so that the action on $J_{x}^{r} B M$ factors through the homomorphic projection $\operatorname{Diff}(M, x) \rightarrow G_{n}^{r}$, where $G_{n}^{r}$ is a finite-dimensional Lie group known as the $r$-jet group $[16, \S 13]$. Thus, we need only consider the invariance of $F_{x}$ under $G_{n}^{r}$. The $r$-jet groups come with natural projections $G_{n}^{r} \rightarrow G_{n}^{r-1}$, corresponding to the equivariant projection $J_{x}^{r} B M \rightarrow J_{x}^{r-1} B M$, and it is easily seen that $G_{n}^{1} \cong$ $G L(n)$. Analogously, for orientation preserving diffeomorphisms, we denote the corresponding projections as Diff ${ }^{+}(M) \rightarrow G_{n}^{+r} \rightarrow G L^{+}(n)$.

The curvature coordinates $\left(g_{a b}, \Gamma_{(b c, A)}^{a}, \bar{S}^{a b(c d, A)}, \bar{z}^{A}\right)$ are used specifically for their transformation properties under $G_{n}^{r}$. Note that, without loss of generality but after a possible small restriction of $V_{x}^{r}$, we can factor $V_{x}^{r} \cong \mathbb{R}^{\gamma} \times W^{r}$, where the projection onto the $\mathbb{R}^{\gamma}$ factor is effected by the $\left(\Gamma_{(b c, A)}^{a}\right)$ coordinates and the projection onto the $W^{r}$ factor is effected by the remaining coordinates. This factorization respects the action of $G_{n}^{r}$ in the sense that the projection $V_{x}^{r} \rightarrow$ $W^{r}$ induces a well-defined action of $G_{n}^{r}$ and $W^{r}$. The action on $W^{r}$ actually factors through the projection $G_{n}^{r} \rightarrow G_{n}^{1} \cong G L(n)$, since it is coordinatized by components of tensors. Moreover, for any $w \in W^{r}$, the isotropy subgroup of $w$ in $G_{n}^{r}$ acts transitively on the fiber $\mathbb{R}^{\gamma}$ over $w$. In the orientation preserving case, the same is true of the corresponding actions of $G_{n}^{+r}$ and $G L^{+}(n)$. The fact that $G_{n}^{r}$ (and also $G_{n}^{+r}$ ) acts transitively on the $\mathbb{R}^{\gamma}$ fibers that are coordinatized by the derivatives of the Christoffel symbols $\left(\Gamma_{(b c, A)}^{a}\right)$ means that an invariant function $F_{x}$ cannot depend on these coordinates, which is a well-known result that is sometimes known as the Thomas replacement theorem [15, 2]. Let us rephrase it slightly below.

The above factorization $V^{r} \cong \mathbb{R}^{\gamma} \times W^{r}$ is also compatible with the rescaled curvature coordinates

$$
\begin{equation*}
\left(x^{a}, g, g^{-\frac{1}{n}} g_{a b}, \Gamma_{(b c, A)}^{a}, g^{\frac{3}{n}+\frac{1}{n}|A|} \bar{S}^{a b(c d, A)}, g^{\frac{s}{2 n}+\frac{1}{n}|A|} \bar{z}^{A}\right) \tag{25}
\end{equation*}
$$

that were introduced in Sect. 2.2. Recall that in our notation the functions $\left(g^{-\frac{1}{n}} g_{a b}\right)$ are functionally independent only up to the identity $\left|\operatorname{det}\left(g^{-\frac{1}{n}} g_{a b}\right)\right|=$ 1. The main distinction is that these coordinates, other than $\left(x^{a}, \Gamma_{(b c, A)}^{a}\right)$, are no longer components of tensors, but rather of tensor densities, which also transform under $G L(n)$ (cf. Sect. 2.6). Using these coordinates, together with the preceding discussion, we can simplify a $\operatorname{Diff}(M)$-invariant $F$ as follows:

Proposition 2.6 (Thomas replacement theorem). Let $F: V_{x}^{\prime r} \subseteq J^{r} B M \rightarrow \mathbb{R}$ be a $\operatorname{Diff}(M)$-invariant function defined on a $\operatorname{Diff}(M)$-invariant domain. In the coordinate system (25) defined on the domain $V^{r} \subseteq V^{\prime r}$, the restriction of $F$ to $V^{r}$ must be expressible as

$$
\begin{equation*}
F=G\left(g, g^{-\frac{1}{n}} g_{a b}, g^{\frac{3}{n}+\frac{1}{n}|A|} \bar{S}^{a b(c d, A)}, g^{\frac{s}{2 n}+\frac{1}{n}|A|} \bar{z}^{A}\right) \tag{26}
\end{equation*}
$$

where the function $G$ is invariant under the action of $G L(n)$ on its arguments.
At this point, we have reduced the invariance of $F$ under $\operatorname{Diff}(M)$ to the invariance of the function $G$, from Proposition 2.6, under $G L(n)$ (obtained as the
projection $\operatorname{Diff}(M, x) \rightarrow G L(n))$, which follows from the preceding discussion. Analogous statements hold for $\mathrm{Diff}^{+}(M)$, $\mathrm{Diff}^{+}(M, x)$ and $G L^{+}(n)$. We now single out a specific subgroup of $G L^{+}(n)$ (and hence also of $G L(n)$ ) that we shall call the group of coordinate scalings. It consists of matrices of the form $\mu I_{n} \in$ $G L(n)$, where $\mu$ is a positive real number and $I_{n}$ is the $n \times n$ identity matrix. The name refers to the fact that $\mu I_{n}$ is the image of a diffeomorphism that restricts to a uniform scaling of the coordinates $\left(x^{a}\right)$ centered at $x \in U \subseteq M$, with of course many other possible pre-images, under the projection $\operatorname{Diff}^{+}(M) \rightarrow G L^{+}(n)$. These transformations should be contrasted with the distinct group of physical scalings introduced in Sect. 2.4.

Coordinate scalings act on the components of tensor densities appearing in the coordinate system (25) as follows:

$$
\begin{align*}
g & \mapsto \mu^{2 n} g, & g^{\frac{3}{n}+\frac{1}{n}|A|} \bar{S}^{a b(c d, A)} & \mapsto \mu^{2+|A|} g^{\frac{3}{n}+\frac{1}{n}|A|} \bar{S}^{a b(c d, A)},  \tag{27}\\
g^{-\frac{1}{n}} g_{a b} & \mapsto g^{-\frac{1}{n}} g_{a b}, & g^{\frac{s}{2 n}+\frac{1}{n}|A|} \bar{z}^{A} & \mapsto \mu^{s+|A|} g^{\frac{s}{2 n}+\frac{1}{n}|A|} \bar{z}^{A} . \tag{28}
\end{align*}
$$

We stress a fundamental difference between coordinate scalings and the previously introduced physical scalings: coordinate scalings are induced from the action of the diffeomorphism group, while the physical ones are not.

### 2.6 Equivariant and isotropic tensors

In this section, we present some basic facts about equivariant maps between spaces that carry certain representations of $G L(n)$.

In particular, consider the space $B_{n}$ of bilinear forms on $\mathbb{R}^{n}$, and the natural linear action of $G L(n)$ thereon. The subset $L_{n} \subset B_{n}$ of non-degenerate bilinear forms of Lorentzian signature $(-+\cdots+)$ is invariant and hence inherits an action of $G L(n)$ itself. If $\eta \in L_{n}$ is the canonical Lorentzian form, defined by the matrix $\operatorname{diag}(-1,1, \ldots, 1)$ referring to the canonical basis of $\mathbb{R}^{n}$, the subgroup $O(1, n-1) \subset G L(n)$ is defined as the isotropy group of $\eta$. We could also restrict the action on $L_{n}$ to the subgroup $G L^{+}(n) \subset G L(n)$ of orientation preserving transformations. With this choice, the isotropy group of $\eta$ turns out to be $S O(1, n-1)=O(1, n-1) \cap G L^{+}(n)$.
Remark 2.2. $L_{n}$ consists of a single orbit and is in fact isomorphic to the homogeneous space $G L(n) / O(1, n-1)$. Similarly, $L_{n}$ is also isomorphic to the homogeneous space $G L^{+}(n) / S O(1, n-1)$. The fact that the action of $G L(n)$ (resp. $\left.G L^{+}(n)\right)$ is transitive on $L_{n}$ implies, as a general well-known fact, that the isotropy group of any $g \in L_{n}$ is isomorphic to $O(1, n-1)$ (resp. $S O(1, n-1)$ ).

Definition 2.5. Let $M_{n}^{p}$ be the space of $p$-multilinear forms on $\mathbb{R}^{n}$ and consider the natural linear action of $G L(n)$ thereon. Let $T$ be a finite-dimensional real vector space carrying a representation of $G L(n)$.
(a) $T$ is a (covariant) tensor representation if it is the restriction of the action of $G L(n)$ on $M_{n}^{p}$ with respect to some linear embedding $T \hookrightarrow M_{n}^{p}$ as an invariant subspace. We call $p$ the tensor rank of $T$.
(b) $T$ a (covariant) tensor density representation if $T$ is as in (a) but the action of $G L(n) \ni u \mapsto \rho(u)$ on $T$ is given by a tensor representation up to a multiplication by $|\operatorname{det} u|^{s}$, where $s$ is the tensor weight of $T$.

Of course, we obtain similar definitions by substituting $G L^{+}(n)$ for $G L(n)$, and also $O(1, n-1)$ or $S O(1, n-1)$, when a particular Lorentzian bilinear form $g$ is fixed. Of course, in the case of $O(1, n-1)$ and $S O(1, n-1)$, there is no distinction between tensor and tensor density representations.

Finally, it is useful to consider the one point space $* \cong \mathbb{R}^{0}$ with the trivial action of $G L(n)$ or any of its subgroups thereon.

Definition 2.6. Let $X$ and $Y$ be spaces carrying respective actions $\rho^{(X)}$ and $\rho^{(Y)}$ of the group $G$. A map $f: X \rightarrow Y$ is said to be equivariant if it commutes with the action of $G$ :

$$
\begin{equation*}
f \circ \rho_{u}^{(X)}=\rho_{u}^{(Y)} \circ f \quad \text { for every } u \in G \tag{29}
\end{equation*}
$$

Consider the special case where $X:=*, Y:=T$ as in (a) in Definition 2.5, and $G:=O(1, n-1)$. The image of an equivariant map $* \rightarrow T$ is called an $O(1, n-1)$-isotropic tensor. The space of $O(1, n-1)$-isotropic tensors in $T$ will be denoted by $\mathcal{I}_{T}$.

An $S O(1, n-1)$-isotropic tensor is defined similarly, replacing $O(1, n-1)$ by $S O(1, n-1)$ everywhere. The space of $S O(1, n-1)$-isotropic tensors in $T$ will be denoted by $\tilde{\mathcal{I}}_{T}$.

## Remark 2.3.

(1) The embedding $T \hookrightarrow M_{n}^{p}$ is an evident example of equivariant map for $G L(n)$ (and every subgroup) by definition.
(2) As $f: * \rightarrow T$ is completely defined by its image $f(*)=t \in T$ the definition states that a tensor $t \in T$ is isotropic if it is invariant under the relevant action of $O(1, n-1)$ (or $S O(1, n-1)$ ) on $T$.
(3) The space of isotropic tensors for different Lorentzian bilinear forms are clearly isomorphic.

It is well known that the subspaces of isotropic tensors $\mathcal{I}_{T} \subset T$ and $\tilde{\mathcal{I}}_{T} \subset T$ can be fully characterized as in the proposition below. In the following, $\epsilon \in M_{n}^{n}$ denotes the canonical Levi-Civita tensor, that is, the fully anti-symmetric form uniquely fixed by the value of its component $\epsilon_{1 \cdots n}=1$, with respect to the canonical basis of $\mathbb{R}^{n}$. Also, $\mathcal{I}_{n}^{p} \subset M_{n}^{p}$ will denote the subspace spanned by all possible tensor products of the canonical Lorentzian form $\eta \in L_{n}$ that create a $p$-multilinear form. More precisely, $\mathcal{I}_{n}^{p}$ is spanned by elements of the form

$$
\begin{equation*}
\left(\eta_{\sigma}\right)_{i_{1} i_{2} \cdots i_{p-1} i_{p}}=\eta_{\sigma\left(i_{1}\right) \sigma\left(i_{2}\right)} \cdots \eta_{\sigma\left(i_{p-1}\right) \sigma\left(i_{p}\right)}, \tag{30}
\end{equation*}
$$

where $\sigma \in S_{p}$ is any permutation. Similarly, $\tilde{\mathcal{I}}_{n}^{p} \subset M_{n}^{p}$ denotes the subspace spanned by all possible tensor products of $\eta$ and $\epsilon$ that create a $p$-multilinear form.

Proposition 2.7. Given a real vector space $T$ carrying a tensor representation of $G L(n)$ and identifying $T$ with its image with respect to the embedding $\alpha$ : $T \hookrightarrow M_{n}^{p}$, the following facts hold.
(a) The subspace $\mathcal{I}_{\tilde{\tilde{I}}} \subset T$ is given by $\mathcal{I}_{\tilde{\mathcal{I}}} \cong \alpha(T) \cap \mathcal{I}_{n}^{p}$.
(b) The subspace $\tilde{\mathcal{I}}_{T} \subset T$ is given by $\tilde{\mathcal{I}}_{T} \cong \alpha(T) \cap \tilde{\mathcal{I}}_{n}^{p}$.

An elementary proof of such a characterization of $O(n)$ - and $S O(n)$-isotropic tensors can be found in [3], which generalizes straightforwardly to $O(1, n-1)$ and $S O(1, n-1)$. More generally, this kind of result is sometimes known as first fundamental theorem of invariant theory [27, 12] for the corresponding group.
Definition 2.7. Given a real vector space $T$ with a tensor density representation of $G L(n)$ (resp. $G L^{+}(n)$ ) and the natural representation on $L_{n}$, we will refer to an equivariant map $t: L_{n} \rightarrow T$ as a $G L(n)$-equivariant tensor density, and similarly for $G L^{+}(n)$-equivariant tensor densities.

The space of $G L(n)$-equivariant tensor densities will be denoted by $\mathcal{E}_{T}$ and the space of $G L^{+}(n)$-equivariant tensor densities will be denoted by $\tilde{\mathcal{E}}_{T}$.
Remark 2.4. Even if the functions belonging to $\mathcal{E}_{T}$ and $\tilde{\mathcal{E}}_{T}$ are not required to be linear, these spaces enjoy a natural structure of real vector space, just in view of the fact that the equivariant tensor densities are maps with values in the real vector space $T$.

The following lemma characterizes the space of equivariant tensor densities (in the sense of equivariant maps) in terms of isotropic tensors (in the sense of the subspaces $\mathcal{I}_{T} \subseteq T$ (resp. $\tilde{\mathcal{I}}_{T} \subseteq T$ ) defined earlier).

Lemma 2.8. Let $T$ be a finite-dimensional real vector space carrying a tensor density representation of $G L(n)$, resp. $G L^{+}(n)$, and assume that $L_{n}$ is equipped with the natural representation.
(a) The space of $G L(n)$-equivariant, resp. $G L^{+}(n)$-equivariant, tensor densities is isomorphic the subspace of $O(1, n-1)$-isotropic tensors, resp. $S O(1, n-$ $1)$-isotropic tensors, in T. More precisely, the isomorphism is defined by

$$
\begin{equation*}
\mathcal{E}_{T} \ni t \mapsto t(\eta) \in \mathcal{I}_{T} \quad\left(\text { resp. } \quad \tilde{\mathcal{E}}_{T} \ni t \mapsto t(\eta) \in \tilde{\mathcal{I}}_{T}\right) \tag{31}
\end{equation*}
$$

(b) For a given $t \in \mathcal{E}_{T}$, we have

$$
\begin{equation*}
t(g)=|\operatorname{det} g|^{s} P(g) \quad \text { for all } \quad g \in L_{n} \tag{32}
\end{equation*}
$$

where $P(g)$ is a homogeneous $T$ valued polynomial in the components of $g$ (with respect to the canonical basis of $\mathbb{R}^{n}$ ), and $s$ is some real number fixed by weight of the tensor density representation of $G L(n)$.
(c) For a given $t \in \tilde{\mathcal{E}}_{T}$, we have

$$
\begin{equation*}
t(g)=|\operatorname{det} g|^{s} P(g, \varepsilon(g)) \quad \text { for all } \quad g \in L_{n}, \tag{33}
\end{equation*}
$$

where $P(g, \varepsilon(g))$ is a homogeneous $T$ valued polynomial in the components of $g$ and the components ${ }^{4}$ of $\varepsilon(g):=\sqrt{\operatorname{det} g} \epsilon$ (with respect to the canonical basis of $\mathbb{R}^{n}$ in both cases), and $s$ is some real number fixed by the weight of the tensor density representation of $G L^{+}(n)$.
Remark 2.5. Since, in view of this Lemma, an equivariant tensor density $t(g)$ is a homogeneous function, say of degree $k$, of $g$ up to a power of $|\operatorname{deg} g|$, it could always be rewritten as

$$
\begin{equation*}
t(g)=|\operatorname{det} g|^{\frac{k}{n}} t\left(|\operatorname{det} g|^{-\frac{1}{n}} g\right) \tag{34}
\end{equation*}
$$

This observation will be later useful in the proof of Theorem 3.1.

[^4]Proof. We deal with the $G L(n)$-equivariant case, the $G L^{+}(n)$-equivariant case being completely analogous. The action of $G L(n)$ on both $L_{n}$ and $T$ is linear, so we denote it as $u \cdot x$, for $u \in G L(n)$ and $x$ in either $L_{n}$ or $T$.

The first crucial observation, as $L_{n}$ consists of a single orbit of $G L(n)$, is that equivariance allows us to fully fix $t: L_{n} \rightarrow T$ provided that we know its value on $\eta \in L_{n}$, by the formula

$$
\begin{equation*}
t(g)=t\left(u_{g} \cdot \eta\right)=u_{g} \cdot t(\eta) \tag{35}
\end{equation*}
$$

for any $g \in L_{n}$ and $u_{g} \in G L(n)$ such that $g=u_{g} \cdot \eta$. The second crucial observation is that, to make sure that the values of $t$ are assigned consistently, $t(\eta)$ must be invariant under the isotropy subgroup of $\eta$, namely $O(1, n-1)$. In other words, $t(\eta)$ must belong to $\mathcal{I}_{T}$, with respect to the induced representation of $O(1, n-1)$ on $T$. The formula (35) clearly defines mutually inverse maps $\mathcal{E}_{T} \rightarrow \mathcal{I}_{T}$ and $\mathcal{I}_{T} \rightarrow \mathcal{E}_{T}$, thus establishing the isomorphism $\mathcal{E}_{T} \cong \mathcal{I}_{T}$ claimed in part (a).

Let us now prove part (b). Fix an (equivariant) embedding $\alpha: T \rightarrow M_{n}^{p}$. Since $t(\eta)$ is an element of $\mathcal{I}_{T}$, from the characterization of isotropic tensors in Proposition 2.7, it must be of the form

$$
\begin{equation*}
t(\eta)=\alpha^{-1}\left(\sum_{\sigma \in S_{p}} c^{\sigma} \eta_{\sigma}\right) \tag{36}
\end{equation*}
$$

where $c^{\sigma}$ are some scalar coefficients. Then for any $g \in L_{n}$ and a corresponding $u_{g} \in G L(n)$ such that $g=u_{g} \cdot \eta$,

$$
\begin{align*}
t(g)=u_{g} \cdot t(\eta) & =\alpha^{-1}\left(\left|\operatorname{det} u_{g}\right|^{r} \sum_{\sigma \in S_{p}} c^{\sigma}\left(u_{g} \cdot \eta_{\sigma}\right)\right) \\
& =|\operatorname{det} g|^{r / 2} \alpha^{-1}\left(\sum_{\sigma \in S_{p}} c^{\sigma} g_{\sigma}\right) \tag{37}
\end{align*}
$$

where $r$ is the density weight of the representation $T$ and where we have used the notation

$$
\begin{equation*}
g_{\sigma}=g_{\sigma\left(i_{1}\right) \sigma\left(i_{2}\right)} \cdots g_{\sigma\left(i_{p-1}\right) \sigma\left(i_{p}\right)} \tag{38}
\end{equation*}
$$

for the corresponding monomial on $L_{n}$ in terms of the components of $g$ with respect to the canonical basis on $\mathbb{R}^{n}$. Clearly, the above formula can be rewritten as $t(g)=|\operatorname{det} g|^{s} P(g)$, with $s:=r / 2$. We observe that, from (37), that $P$ is an homogeneous polynomial (of degree $p / 2$ ) in the components of the metric, completing the proof of part (b).

The proof of (c) is strictly analogous, taking into account the identity $u_{g} \cdot \epsilon=$ $\varepsilon(g)$, for any $u_{g} \in G L^{+}(n)$ such that, $u_{g} \cdot \eta=g$.

## 3 Characterization of Finite Renormalizations of Wick Polynomials

We generalize the discussion of local covariant fields from [13], where only metric dependence was allowed, to a more general context where other background
fields are allowed in addition to the metric $\mathbf{g}$ on a spacetime $M$. In order to simplify the presentation, we will restrict the extra background fields to two scalar functions $m^{2}$ and $\xi$, which appear in the description of a scalar quantum field.

Generally speaking, background fields are described by sections $\mathbf{h}$ of suitable bundles $H M \rightarrow M$ over the manifolds $M$ we consider. Covariance requires us to deal with all such bundles simultaneously and coherently. In other words we deal with an assignment of a bundle $H M \rightarrow M$ to every manifold $M$ and require that any embedding $\chi: M \rightarrow M^{\prime}$ must give rise to a corresponding well-defined pullback map $\chi^{*}: \Gamma\left(H M^{\prime}\right) \rightarrow \Gamma(H M)$. This picture can be phrased properly with the language of category theory by means of the notion of natural bundle. Before giving the definition, we would like to require a bit more geometric structure from the bundles of background fields that interest us. A bundle $F \rightarrow M$ is dimensionful if there it has an action $\mathbb{R}^{+} \times F \rightarrow F$ of the multiplicative group $\mathbb{R}^{+}$of positive real numbers, called (physical) scaling, which acts by a diffeomorphisms that fix each fiber of $F \rightarrow M$. Any vector bundle is automatically dimensionful, by virtue of having a well-defined multiplication by scalars on its fibers, although we will not always use this particular scaling action. To avoid confusion, we should mention that a dimensionless bundle would be a special kind of dimensionful bundle, where scaling transformations act trivially.

A natural (dimensionful) bundle is a functor $H: \mathfrak{M a n} \rightarrow \mathfrak{B n d l}$ from the category of smooth manifolds (where objects are connected, have fixed dimension $n$ and morphisms are embeddings, which are necessarily local diffeomorphism) to the category of dimensionful smooth bundles (where morphisms are bundle maps, i.e., fiber preserving, equivariant with respect to scaling), such that a morphism $\chi: M \rightarrow M^{\prime}$ induces a morphism $H \chi: H M \rightarrow H M^{\prime}$ that is itself a local diffeomorphism. The required pullback $\chi^{*}: \Gamma\left(H M^{\prime}\right) \rightarrow \Gamma(H M)$ is then implicitly defined by $\mathbf{h}^{\prime} \circ \chi=H \chi \circ\left(\chi^{*} \mathbf{h}^{\prime}\right)$, when $\mathbf{h}^{\prime} \in \Gamma\left(H M^{\prime}\right)$. The equivariance of the morphism $H \chi$ ensures that scaling commutes with the pullback, $\chi^{*}\left(\mathbf{h}_{\lambda}^{\prime}\right)=\left(\chi^{*} \mathbf{h}^{\prime}\right)_{\lambda}$, for $\lambda \in \mathbb{R}^{+}$.

One elementary example of a natural bundle is the functor $M \mapsto \mathbb{R} \times M$, the trivial scalar bundle, whose sections we call scalar fields, with scaling being simple multiplication. Another relevant example is $M \mapsto S^{2} T^{*} M$, the bundle of Lorentzian metrics; we will denote a section of $\dot{S}^{2} T^{*} M \rightarrow M$ by g. Other examples are are $M \mapsto T^{*} M$ and $M \mapsto \Lambda^{2} M$, the cotangent bundle and the bundle of 2 -forms, whose sections could be interpreted as background electromagnetic fields, in the vector potential or field strength forms. All of these bundles are dimensionful, by virtue of being vector bundles, with the exception of $\stackrel{\circ}{S}^{2} T^{*} M$, which inherits a scaling action from being considered as a scaling invariant sub-bundle of the vector bundle $S^{2} T^{*} M$.

Remark 3.1. In the rest of the paper, focussing on the theory of a real quantum scalar field, $\varphi$, we make a more precise choice of the natural functor $H$. We suppose that the manifolds of the category $\mathfrak{M a n}$ are connected, $n$-dimensional (for a fixed $n \geq 2$ ), and the functor $H$ assigns $M \mapsto H M=S^{2} T^{*} M \times \mathbb{R} \times \mathbb{R}$, with a morphism $\chi: M \rightarrow M^{\prime}$ inducing the standard tensor push-forward $H \chi=$ $\chi_{*}: H M \rightarrow H M^{\prime}$. Then, the sections $M \rightarrow H M$ are triples $\mathbf{h}=\left(\mathbf{g}, m^{2}, \xi\right)$, scaling as $\left(\mathbf{g}, m^{2}, \xi\right) \mapsto\left(\lambda^{-2} \mathbf{g}, \lambda^{2} m^{2}, \xi\right)$, always consisting of:
(a) a Lorentzian metric, $\mathbf{g}$, making $(M, \mathbf{g})$ a (smooth) Lorentzian spacetime of fixed dimension $n \geq 2$,
(b) the pair of real scalar fields $m^{2}$ and $\xi$ over $M$, with the respective physical meaning of the squared mass of the scalar field and a factor describing the coupling with the scalar curvature.

We stress that, exactly as in [13], we assume that the parameters $m^{2}$ and $\xi$ are actually functions on $M$. Quantum field theory in curved spacetime is welldefined for both constant or variable $m^{2}$ and $\xi$. There is of course no obstacle in restricting them to constant functions, as we note in Remark 3.4. Moreover, as in [13], $m^{2}$ and $\xi$ are allowed to have any real value.

Definition 3.1. Let us fix the natural bundle $H: \mathfrak{M a n} \rightarrow \mathfrak{B n d}$ as in Remark 3.1. A background field is a section $\mathbf{h}: M \rightarrow H M$ and we call the pair $(M, \mathbf{h})$ a background geometry, provided $\mathbf{h}=\left(\mathbf{g}, m^{2}, \xi\right)$ is such that $(M, \mathbf{g})$ is a time-orientable globally hyperbolic spacetime. Furthermore we define the following categories.
(a) $\mathfrak{B k g} \mathfrak{G}$ is the category of background geometries, having timeoriented background geometries as objects and morphisms given by smooth embeddings $\chi: M \rightarrow M^{\prime}$ that preserve the background fields, $\chi^{*} \mathbf{h}=\mathbf{h}^{\prime}$ on $M^{\prime}$, the time orientation, and causality, meaning that every causal curve between $\chi(p)$ and $\chi(q)$ in $M^{\prime}$ is the $\chi$-image of a causal curve between $p$ and $q$ in $M$.
(b) $\mathfrak{B k g}^{+}$is the category of oriented background geometries having oriented and time-oriented background geometries as objects and morphisms as in $\mathfrak{B e g} \mathfrak{G}$, but also required to preserve the spacetime orientation.

Since the natural bundle $H$ is dimensionful, scaling transformations also act on these categories by $(M, \mathbf{h}) \mapsto\left(M, \mathbf{h}_{\lambda}\right)$, for any $\lambda \in \mathbb{R}^{+}$, which by equivariance of the pullback of background fields act as functors, $\mathfrak{B k g} \mathfrak{G} \rightarrow \mathfrak{B k g} \mathfrak{G}$ and $\mathfrak{B k g G ^ { + }} \rightarrow \mathfrak{B k g G}^{+}$respectively.

To describe the algebras of observables on background geometries, we need the notion of a net of algebras (or pre-cosheaf of algebras).

Definition 3.2. A net of algebras (of observables) is an assignment of a complex unital $*$-algebra $\mathcal{W}(M, \mathbf{h})$ for every background geometry $(M, \mathbf{h})$ in $\mathfrak{B k g} \mathfrak{G}$ together with an assignment of an injective unital $*$-algebra homomorphism $\iota_{\chi}: \mathcal{W}(M, \mathbf{h}) \rightarrow \mathcal{W}\left(M^{\prime}, \mathbf{h}^{\prime}\right)$ for every morphism in $\mathfrak{B k g} \mathfrak{G}$, respecting compositions. In other words $\mathcal{W}: \mathfrak{B k g} \mathfrak{G} \rightarrow \mathfrak{A l g}$ is a functor from the category of background geometries into the category of (complex) unital $*$-algebras whose morphisms are injective unital $*$-algebra homomorphisms. Further, we require that $\mathcal{W}$ respects scaling and the time slice axiom.
(i) Scaling transformations $(M, \mathbf{h}) \mapsto\left(M, \mathbf{h}_{\lambda}\right)$ result in $*$-algebra isomorphisms $\sigma_{\lambda}: \mathcal{W}(M, \mathbf{h}) \rightarrow \mathcal{W}\left(M, \mathbf{h}_{\lambda}\right)$. Scaling transformations act as natural isomorphisms $\sigma_{\lambda}: \mathcal{W} \rightarrow \mathcal{W}_{\lambda}$ between the $*$-algebra valued functors $\mathcal{W}$ and $\mathcal{W}_{\lambda}$, the latter defined by $\mathcal{W}_{\lambda}(M, \mathbf{h})=\mathcal{W}\left(M, \mathbf{h}_{\lambda}\right)$.
(ii) Given a morphism $\chi:\left(M^{\prime}, \mathbf{h}^{\prime}\right) \rightarrow(M, \mathbf{h})$ of background geometries, if the image $\chi\left(M^{\prime}\right) \subseteq M$ contains a Cauchy surface for $(M, \mathbf{g})$, then the induced *-homomorphism $\iota_{\chi}: \mathcal{W}\left(M^{\prime}, \mathbf{h}^{\prime}\right) \rightarrow \mathcal{W}(M, \mathbf{h})$ is a $*$-isomorphism.

We refer to a functor $\mathcal{W}: \mathfrak{B k g G}^{+} \rightarrow \mathfrak{A l g}$ with the same properties as a net of algebras as well.

The algebras of observables $\mathcal{W}(M, \mathbf{h})$ are intended to be the algebras of Wick products, whose construction and the fact that they satisfy all the desired properties are discussed in detail in [13, Sec. 2], cf. their Lemma 4.2 in particular, where they define the scaling isomorphism $\sigma_{\lambda}$. However, we shall not touch upon these details and rely only on the properties of the $\mathcal{W}$ functor as axiomatized above and also below in the definition of Wick powers.

Having defined a net of algebras, respecting local covariance and scaling, we are in a position to state our definition of a locally covariant (almost) homogeneous quantum field which somewhat extends and generalizes [13, Def. 3.2].
Definition 3.3. A locally covariant scalar quantum field $\Phi$ is an assignment of an algebra-valued distribution ${ }^{5} \Phi_{(M, \mathbf{h})}: \mathcal{D}(M) \rightarrow \mathcal{W}(M, \mathbf{h})$ to each background geometry $(M, \mathbf{h})$ that satisfies the following identity for each morphism $\chi:\left(M^{\prime}, \mathbf{h}^{\prime}=\chi^{*} \mathbf{h}\right) \rightarrow(M, \mathbf{h}):$

$$
\begin{equation*}
\iota_{\chi}\left(\Phi_{\left(M, \chi^{*} \mathbf{h}\right)}(f)\right)=\Phi_{(M, \mathbf{h})}\left(\chi_{*} f\right), \quad \text { for any } f \in \mathcal{D}\left(M^{\prime}\right) \tag{39}
\end{equation*}
$$

In other words, $\Phi$ is a natural transformation $\Phi: \mathcal{D} \rightarrow \mathcal{W}$ between the functor of test functions ${ }^{6}$ and the net of algebras of observables (thus, another name for $\Phi$ could be a natural scalar quantum field).

The reason that a quantum field $\Phi$, even on a fixed spacetime $(M, \mathbf{g})$, is associated with an algebra-valued distributions is the usual heuristic according to which pointlike fields $\Phi(x)$ are too singular to be evaluated directly, while its smearing with a test function $f \in \mathcal{D}(M)$,

$$
\begin{equation*}
\Phi(f)=\int_{M} \Phi(x) f(x) d g(x) \tag{40}
\end{equation*}
$$

where $d g(x)$ is the volume form induced by the metric $\mathbf{g}$, is a legitimate observable. Where appropriate, we will use the distributional notation $\Phi(x)$ as well.

Note the most trivial example of a locally covariant scalar quantum field, which we may call the unit c-number field $\mathbf{1}$, defined by the formula

$$
\begin{equation*}
\mathbf{1}_{(M, \mathbf{h})}(f):=1 \int_{M} f(x) d g(x) \tag{41}
\end{equation*}
$$

where $1 \in \mathcal{W}(M, \mathbf{h})$ is the algebra unit. If $C[\mathbf{h}]$ is any function that maps a background geometry on $M$ to a distribution on $M$ and satisfies the identity $\chi^{*} C[\mathbf{h}]=C\left[\chi^{*} \mathbf{h}\right]$ for any morphism $\chi:\left(M^{\prime}, \chi^{*} \mathbf{h}\right) \rightarrow(M, \mathbf{h})$, then the product $C \mathbf{1}$ functor defined by the formula

$$
\begin{equation*}
C \mathbf{1}_{(M, \mathbf{h})}(f)=1 \int_{M} C[\mathbf{h}](x) f(x) d g(x) \tag{42}
\end{equation*}
$$

is also a locally covariant scalar quantum field that we refer to as a c-number field.

[^5]Remark 3.2.
(1) In [13], $\mathbf{h}$ is nothing but the Lorentzian metric of the spacetime and the parameters $m^{2}$ and $\xi$ appearing in the definition of the quantum fields generated by KG fields are considered external parameters. Here instead we explicitly include them in $\mathbf{h}$. It is very easy to prove that the concrete locally covariant quantum fields appearing in [13] (scalar KG field and associated Wick polynomials, time-ordered Wick polynomials and their derivatives) satisfy our somewhat more general definition locally covariant quantum fields.
(2) Definition 3.2 includes two distinct though related notions: locality and covariance, both illustrated by the condition (39). Locality corresponds to the case where $\chi$ describes the inclusion $\chi: M \subset M^{\prime}$, while covariance corresponds to an arbitrary allowed $\chi$.

In [13, Sec. 2], the algebras $\mathcal{W}(M, \mathbf{h})$ are constructed explicitly for the case of the quantization of a Klein-Gordon scalar field $\varphi$, with mass $m^{2}$ and coupled with the scalar curvature through the constant $\xi$, whose equation of motion is

$$
\begin{equation*}
\square_{\mathbf{g}} \varphi-m^{2} \varphi-\xi R_{\mathbf{g}} \varphi=0 \tag{43}
\end{equation*}
$$

In that context, the basic example of a locally covariant scalar field is the KleinGordon field $\varphi$ itself. Since we are not dealing with such explicit constructions, we simply encapsulate the needed properties of $\varphi$ in the definition below.

Definition 3.4. Given a net of algebras $\mathcal{W}$ on $\mathfrak{B k g G}$ or $\mathfrak{B k g} \mathfrak{G}^{+}$, a linear quantum scalar field $\varphi$ is a locally covariant quantum scalar field that satisfies the following kinematic completeness property: For any ( $M, \mathbf{h}$ ), an element $a \in \mathcal{W}(M, \mathbf{h})$ satisfies $\left[a, \varphi_{(M, \mathbf{h})}(f)\right]=0$ for every $f \in \mathcal{D}(M)$ iff $a=\alpha 1$, with $\alpha \in \mathbb{C}$ and 1 the unit element of the algebra.

The above definition is rather minimal and, certainly, the linear KG field as defined in [13, Sec. 2] satisfies several additional properties, but we have omitted most of them. Below, we will define Wick powers $\varphi^{k}$, which will include the linear field $\varphi^{1}:=\varphi$ as a special case. Within the definition to follow, we will require further properties to hold for $\varphi^{k}$ for each $k$, including $k=1$, thus imposing further axioms also on the linear field and bringing our axiomatization closer to that of [13]. Thus, to avoid some repetition, we find it more economical to state these axioms in a way that is uniform in $k$. As it is the main goal of this work to remove it, the analyticity axiom will not appear below and the continuity axiom will be suitably modified to compensate. On the other hand, we will not make use of the on-shell condition, $\varphi\left(\left(\square_{\mathbf{g}}-m^{2}-\xi R_{\mathbf{g}}\right) f\right)=0$, though the equations of motion appear implicitly through the time slice axiom in Definition 3.2. We have excluded it because it plays no role in our analysis below, since it is restricted to simple Wick powers. However, the on-shell condition will have to be taken into account when expanding the analysis to more general Wick and time-ordered products that involve derivatives of $\varphi$.

Before giving a precise axiomatic definition of Wick powers, we need to address the technical question of how physical scalings and continuous variations of the background geometry can be made to act on locally covariant scalar quantum fields.

First, we address scalings. Given a locally covariant scalar quantum field $\Phi$, we can define a new rescaled locally covariant scalar quantum field $S_{\lambda} \Phi$ by the formula

$$
\begin{equation*}
\left(S_{\lambda} \Phi\right)_{(M, \mathbf{h})}(f)=\sigma_{\lambda}^{-1}\left(\Phi_{\left(M, \mathbf{h}_{\lambda}\right)}\left(\lambda^{n} f\right)\right) \tag{44}
\end{equation*}
$$

where $\lambda \in \mathbb{R}^{+}, n=\operatorname{dim} M$, and $\sigma_{\lambda}$ is the isomorphism realizing the action of scalings on the net of algebras of observables. The extra factor of $\lambda^{n}$ compensates for the fact that integration against the test function $f$ is done with respect to the metric volume form $d g$, which scales as $d g \mapsto \lambda^{-n} d g$ when $\mathbf{g} \mapsto \lambda^{-2} \mathbf{g}$. Comparing this formula with [13, Eq. (48)], note that the direction of our isomorphism $\sigma_{\lambda}$ is opposite. More formally, recall that $\sigma_{\lambda}: \mathcal{W} \rightarrow \mathcal{W}_{\lambda}$ is a natural isomorphism between two functors (defined on $\mathfrak{B k g} \mathfrak{G}$ or $\mathfrak{B k g} \mathfrak{G}^{+}$), while $\Phi: \mathcal{D} \rightarrow \mathcal{W}$ is another natural transformation. If we similarly define the natural linear isomorphism $\mu_{\lambda}: \mathcal{D} \rightarrow \mathcal{D}$ given by $f \mapsto \lambda^{n} f$, the rescaled quantum field is given by a composition of these natural transformations, $S_{\lambda} \Phi=\sigma_{\lambda}^{-1} \circ \Phi \circ \mu_{\lambda}$. As such, we have defined a representation $S_{\lambda}$ of the multiplicative group $\mathbb{R}^{+}$ on the vector space ${ }^{7}$ of natural transformations $\mathcal{D} \rightarrow \mathcal{W}$, which opens the door to using the generic notion of homogeneous and almost homogeneous elements from Definition 2.3 to quantum fields.

Next, in order to address the continuity hypothesis, we must specify how to identify the algebras of observables defined on different background geometries, as was done in [13, Sec. 4.2]. Let $\left(M, \mathbf{h}_{1}\right)$ and $\left(M, \mathbf{h}_{2}\right)$ be two background geometries that differ only inside a compact set $O \subset M$, recalling that $\mathbf{h}_{i}=$ $\left(\mathbf{g}_{i}, m_{i}^{2}, \xi_{i}\right)$. It is a simple fact of Lorentzian geometry that the spacetimes $\left(M_{ \pm}, \mathbf{g}_{ \pm}\right)$, with $M_{-}=M \backslash J^{+}(O), M_{+}=M \backslash J^{-}(O)$ and $\mathbf{h}_{ \pm}=\left.\mathbf{h}_{1}\right|_{M_{ \pm}}=$ $\left.\mathbf{h}_{2}\right|_{M_{ \pm}}$, are globally hyperbolic spacetimes in their own right. Moreover, each of the $\left(M_{ \pm}, \mathbf{h}_{ \pm}\right)$contains a Cauchy surface in common with both $\left(M, \mathbf{h}_{1}\right)$ and $\left(M, \mathbf{h}_{2}\right)$. Thus, denoting by $\chi_{ \pm}^{i}:\left(M_{ \pm}, \mathbf{h}_{ \pm}\right) \rightarrow\left(M, \mathbf{h}_{i}\right)$ the inclusion morphisms, the time slice axiom of the net of the algebras observables gives us isomorphisms ${ }^{{ }_{\chi}{ }_{ \pm}^{i}}: \mathcal{W}\left(M_{ \pm}, \mathbf{h}_{ \pm}\right) \rightarrow \mathcal{W}\left(M, \mathbf{h}_{i}\right)$. These isomorphisms allow us to identify the algebras of observables of background geometries $\left(M, \mathbf{h}_{1}\right)$ and $\left(M, \mathbf{h}_{2}\right)$ that differ only inside a compact $O \subset M$ in two ways, $\tau_{\text {ret }}, \tau_{\text {adv }}: \mathcal{W}\left(M, \mathbf{h}_{1}\right) \rightarrow \mathcal{W}\left(M, \mathbf{h}_{2}\right)$, where these $*$-algebra isomorphisms are defined by $\tau_{\text {ret }}=\iota_{\chi_{2}^{-}} \circ \iota_{\chi_{1}^{-}}^{-1}$ and $\tau_{\text {adv }}=$ ${ }^{\iota} \chi_{2}^{+} \circ \iota_{\chi_{1}^{+}}^{-1}$. Below, we will only make use of $\tau_{\text {ret }}$, though choosing $\tau_{\text {adv }}$ would not have lead to equivalent results.

Finally, we introduce an axiomatic definition of Wick powers, which will be our main objects of interest. We deviate somewhat in our axiomatization from the analogous one in $[13$, Sec. 4] for reasons expanded on below.

Definition 3.5 (Wick powers). Given a net of algebras $\mathcal{W}$ (Definition 3.2) on the category of background geometries $\mathfrak{B k g} \mathfrak{G}$ (or $\mathfrak{B k g} \mathfrak{G}^{+}$) and a corresponding locally covariant linear scalar quantum field $\varphi$ (Definition 3.4), we define Wick powers $\left\{\varphi^{k}\right\}$ of $\varphi$, for $k=0,1,2, \cdots$, by the following axioms:
(i) Locality and Covariance. Each Wick power $\varphi^{k}$ is a locally covariant scalar quantum field (Definition 3.3), which for low powers agree with $\varphi^{0}=1$, the unit $c$-number field, and $\varphi^{1}=\varphi$, the linear field.

[^6](ii) Scaling. Each Wick power $\varphi^{k}$ is almost homogeneous of degree $k(n-2) / 2$ (Definition 2.3) with respect to the action of physical scalings $S_{\lambda}$ (Eq. (44)) on locally covariant fields. That is, there exists an integer $l \geq 0$ and locally covariant fields $\psi_{j}$ such that $S_{\lambda} \varphi^{k}=\lambda^{k \frac{(n-2)}{2}} \varphi^{k}+\lambda^{k \frac{(n-2)}{2}} \sum_{j=1}^{l}\left(\log ^{j} \lambda\right) \psi_{j}$, where each $\psi_{j}$ is itself almost homogeneous of degree $k(n-2) / 2$ and order $l-j$ (as in Eq. (11)). In dimension $n=4,(n-2) / 2=1$ gives the usual scaling of a scalar field.
(iii) Algebraic. Each Wick power $\varphi^{k}$ also satisfies the following properties:

## Hermiticity:

$$
\begin{equation*}
\varphi_{(M, \mathbf{h})}^{k}(f)^{*}=\varphi_{(M, \mathbf{h})}^{k}\left(f^{*}\right) \tag{45}
\end{equation*}
$$

where on the left $*$ denotes the corresponding operation in the $*-$ algebra $\mathcal{W}$, while on the right $*$ denotes simple complex conjugation.

## Commutator expansion:

$$
\begin{equation*}
\left[\varphi_{(M, \mathbf{h})}^{k}(x), \varphi_{(M, \mathbf{h})}(y)\right]=i k \varphi^{k-1}(x) \Delta_{(M, \mathbf{h})}(x, y) \tag{46}
\end{equation*}
$$

where $\Delta_{(M, \mathbf{h})}(x, y)=G_{(M, \mathbf{h})}^{+}(x, y)-G_{(M, \mathbf{h})}^{-}(x, y)$ is the difference between the retarded and advanced Green functions for the KG equation (43).
(iv) Parametrized microlocal spectrum condition. Consider any background geometry $(M, \mathbf{h})$ and any smooth $m$-parameter family $\left(M, \mathbf{H}_{s}\right)$ of compactly supported variations thereof $(m \geq 0)$. That is, $\mathbf{h}=\mathbf{H}_{0}$, where $\mathbf{H}_{s}(x)=\mathbf{H}(s, x)$ with smooth $\mathbf{H}: \mathbb{R}^{m} \times M \rightarrow H M$, such that $\mathbf{H}_{s}$, for any $s \in \mathbb{R}^{m}$, differs from $\mathbf{h}$ only on a compact subset $O \subset M$. Further, let us implicitly identify each algebra of observables $\mathcal{W}\left(M, \mathbf{H}_{s}\right)$ with $\mathcal{W}(M, \mathbf{h})$ using the isomorphism $\tau_{\text {ret }}^{s}: \mathcal{W}\left(M, \mathbf{H}_{s}\right) \rightarrow \mathcal{W}(M, \mathbf{h})$ discussed earlier.
Then, for any quasi-free Hadamard state $\omega$ on $\mathcal{W}(M, \mathbf{h})$ with respect to $\varphi$ and each Wick power $\varphi^{k}$, the expectation value $\omega\left(\varphi^{k}(x)\right)$ is a distribution on $\mathbb{R}^{m} \times M$ with empty wavefront set, which hence can be represented by a smooth function.

Remark 3.3. The above axioms for Wick powers differ in some aspects from those given in [13, Sec. 4].
(a) Our scaling condition, which uses Definition 2.3, is slightly weaker than Definition 4.2 of [13], but it will be sufficient for our purposes. The difference is in the notion of order of the logarithmic terms. The 'order' in Definition 2.3 refers only to the scaling properties. On the other hand, the 'order' of a quantum field used in [13, Def. 4.2] refers to the number of iterated commutations with $\varphi$ needed to annihilate that field. An inductive argument (in $k$ ) shows that if a Wick power $\varphi^{k}$ satisfied Definition 4.2 of [13], then the same Wick power satisfies also our Definition 2.3(b).
(b) The technical continuity and analyticity conditions were replaced by a strengthened version of the microlocal scaling condition. Using the notation from the last point of Definition 3.5, the continuity condition in [13] required $\tau_{\text {ret }}^{s} \circ \varphi_{\left(M, \mathbf{H}_{s}\right)}^{k}(f) \in \mathcal{W}(M, \mathbf{h})$ to be continuous in $s$ for any test function $f \in \mathcal{D}(M)$ and any 1-parameter family $\mathbf{H}(s, x)$, with the topology
on the algebra $\mathcal{W}(M, \mathbf{h})$ left implicit. The analyticity condition required the expectation values $\omega^{s}\left(\varphi_{\left(M, \mathbf{H}_{s}\right)}^{k}(x)\right)$ to be analytic in $(s, x)$ whenever the 1-parameter family $\mathbf{H}(s, x)$ is analytic (of course relaxing the requirement that the variations of the background geometry must vanish outside a compact set) and is accompanied by a 1-parameter family of states $\omega^{s}$ analytic in $s$, with $\omega^{s}$ quasi-free and Hadamard on $\mathcal{W}(M, \mathbf{h})$ and both the linear structure and topology on the appropriate space of states left implicit.
(c) Removing the analyticity condition is our main goal, so that change should not be surprising. In the proof of Theorem 3.1 we will appeal to the Peetre-Slovák theorem, instead of analyticity as was done in [13], to conclude that ambiguities in the Wick powers are characterized by differential operators. Thus, the remaining technical continuity hypothesis is needed only in as much as it helps meet the weak regularity hypothesis needed for the application of the Peetre-Slovák theorem, which is the only one not already covered by locality. However, using the notation from the last point of Definition 3.5, the continuity of $\varphi_{\left(M, \mathbf{H}_{s}\right)}^{k}(f) \in \mathcal{W}(M, \mathbf{h})$ as a function of $s$, even with an opportune choice of topology on $\mathcal{W}(M, \mathbf{h})$, is not sufficient to assure the needed weak regularity property of $\varphi^{k}$. What would be needed instead is a technical infinite-dimensional smoothness condition. Instead of going down that road, we simply strengthen the original microlocal spectrum condition to cover parametrized families of, instead of just individual, background geometries. This new parametrized microlocal condition hypothesis then essentially directly yields the needed weak regularity property of $\varphi^{k}$.

Along with the explicit construction of the algebras $\mathcal{W}(M, \mathbf{h})$ in [13], Hollands and Wald also shows the existence of a family of locally covariant scalar quantum fields that satisfy their version of the axioms of Wick powers, whose difference from ours are discussed above. However, that given construction is not the only way to satisfy these axioms. The lack of uniqueness is physically interpreted as the existence of some remaining degrees of freedom in the renormalization procedure of Wick powers. The key result of [13] on finite renormalizations of Wick powers is stated in [13, Thm. 5.1]. Any pair of families of Wick powers ${ }^{8}\left\{\tilde{\varphi}^{k}\right\}$ and $\left\{\varphi^{k}\right\}(k \in \mathbb{N})$ of the same Klein-Gordon field $\varphi=\varphi^{1}=\tilde{\varphi}^{1}$ satisfies the following relation, in our notation, for every fixed background geometry ( $M, \mathbf{h}$ ):

$$
\begin{equation*}
\tilde{\varphi}_{(M, \mathbf{h})}^{k}(x)=\varphi_{(M, \mathbf{h})}^{k}(x)+\sum_{i=0}^{k-2}\binom{k}{i} C_{k-i}[\mathbf{h}](x) \varphi_{(M, \mathbf{h})}^{i}(x) \tag{47}
\end{equation*}
$$

Above, the scalar coefficients $C_{k}[\mathbf{h}]$ are some scalar differential operators that are tensorially constructed out of the metric, the curvature tensor and its derivatives. These operators depend polynomially on the curvature tensor, its derivatives and on $m^{2}$, with coefficients that depend analytically on $\xi$. Moreover, the $C_{k}[\mathbf{h}]$ scale as $C_{k} \mapsto \lambda^{k \frac{(n-2)}{2}} C_{k}$ when their arguments are rescaled as $\xi \mapsto \xi$, $m^{2} \mapsto \lambda^{2} m^{2}, \mathbf{g}^{a b} \mapsto \lambda^{2} \mathbf{g}^{a b}$ and $R_{a b c d} \mapsto \lambda^{-2} R_{a b c d}$, with the same scaling weight

[^7]for its derivatives. In dimension $n=4$, the coefficients scale as $C_{k} \mapsto \lambda^{k} C_{k}$, which is the expression that appears in [13].

As mentioned in the Introduction, the Analyticity requirement that distinguishes the Hollands and Wald definition of Wick powers from ours (Remark 3.3) is somewhat unnatural and technically difficult to handle, as was stressed in the Introduction. We would like to demonstrate that Analyticity is not necessary to prove a result that is essentially similar to the statement of [13, Thm. 5.1] mentioned above. On the other hand, the Continuity requirement of Hollands and Wald cannot be completely dispensed with, since some version of it must survive to feed into the weak regularity hypothesis needed by the Peetre-Slovák theorem. However, it seems difficult to find a simple modification of Continuity that would do the job. Instead, we find it more useful to strengthen the usual Microlocal Spectrum Condition to its Parametrized version (Definition 3.5(iv)). We believe that this strengthened version is rather natural, encapsulating the stability of the properties of Wick powers under variations of the parameters of the background geometry in precise technical terms, without leaving the bounds of smooth differential geometry.

In [13, Sec. 5.2], Hollands and Wald argue that the by now standard locally covariant Hadamard parametrix prescription for defining Wick powers satisfies all the requirements that we listed in Definition 3.5, with our Parametrized Microlocal Spectral condition replaced by the standard one and adding also their Continuous and Analytic dependence requirements (Remark 3.3). It should be noted that they left the arguments supporting Continuous and Analytic dependence implicit, not giving a complete proof, which appeared later in [14]. Similarly, we believe that the locally covariant Hadamard parametrix prescription satisfies also our Parametrized Microlocal Spectral condition, but leave a detailed complete proof of this claim to future investigations. Thus, by eliminating both the Continuity and Analyticity requirements in favor of the Parametrized Microlocal Spectrum Condition, we can achieve essentially the same result written below into a more precise form:
Theorem 3.1. Let $\left\{\tilde{\varphi}^{k}\right\}$ and $\left\{\varphi^{k}\right\}$ be two families $(k \in \mathbb{N})$ of Wick powers, as in Definition 3.5, with respect to a linear scalar quantum field $\varphi$ (Definition 3.4) in a net of algebras $\mathcal{W}$ (Definition 3.2, defined on either the category $\mathfrak{B k g \mathfrak { G }}$ of nonoriented background geometries or $\mathfrak{B k g G}{ }^{+}$of oriented background geometries).
(a) If $\varphi$ is defined with respect to the category $\mathfrak{B k g} \mathfrak{G}$, for every $(M, \mathbf{h})$, the difference between $\tilde{\varphi}^{k}$ and $\varphi^{k}$ can be parametrized as in (47), where the coefficients

$$
\begin{align*}
C_{k}[\mathbf{h}](x)=C_{k}[ & \mathbf{g}^{a b}(x), R_{a b c d}(x), \ldots, \nabla_{e_{1}} \cdots \nabla_{e_{h}} R_{a b c d}(x), \\
& \left.\xi(x), \ldots, \nabla_{e_{1}} \cdots \nabla_{e_{r}} \xi(x), m(x)^{2}, \ldots, \nabla_{e_{1}} \cdots \nabla_{e_{s}} m(x)^{2}\right] \tag{48}
\end{align*}
$$

are some scalar polynomials, tensorially formed from all of their arguments, except $\xi(x)$, and where $R_{a b c d}(x)$ denotes the Riemann tensor and $\nabla_{a}$ the LeviCivita connection of $\mathbf{g}_{a b}$ at $x \in M$.
(b) If $\varphi$ is defined with respect to the category $\mathfrak{B k g} \mathfrak{G}^{+}$, for every $(M, \mathbf{h})$, we have a variant of (47) with

$$
\begin{align*}
& C_{k}[\mathbf{h}](x)=C_{k}\left[\mathbf{g}^{a b}(x), \varepsilon^{a_{1} \cdots a_{n}}(x), R_{a b c d}(x), \ldots, \nabla_{e_{1}} \cdots \nabla_{e_{h}} R_{a b c d}(x)\right. \\
&\left.\xi(x), \ldots, \nabla_{e_{1}} \cdots \nabla_{e_{r}} \xi(x), m(x)^{2}, \ldots, \nabla_{e_{1}} \cdots \nabla_{e_{s}} m(x)^{2}\right] \tag{49}
\end{align*}
$$

scalar polynomials, tensorially formed from all of their arguments, except $\xi(x)$, and now including the Levi-Civita tensor $\varepsilon^{a_{1} \cdots a_{n}}(x)$ of $\mathbf{g}_{a b}$ at $x \in M$.

In both cases (a) and (b), the coefficients of the polynomials are smooth (instead of analytic) functions of $\xi(x)$ whose functional form does not depend on $M$.

Further, the $C_{k}$ scale as $C_{k} \mapsto \lambda^{k \frac{(n-2)}{2}} C_{k}$ when their arguments are rescaled as follows: $\xi \mapsto \xi, m^{2} \mapsto \lambda^{2} m^{2}, g^{a b} \mapsto \lambda^{2} g^{a b}, \varepsilon^{a_{1} \cdots a_{n}} \mapsto \lambda^{n} \varepsilon^{a_{1} \cdots a_{n}}, R_{a b c d}(x) \mapsto$ $\lambda^{-2} R_{a b c d}(x)$ and the covariant derivatives do not change this rescaling behaviour as the coordinates are dimensionless. These rescaling properties fix the order of the polynomial $C_{k}$.

Obviously all terms $\nabla_{e_{1}} \cdots \nabla_{e_{s}} \xi(x)$ and $\nabla_{e_{1}} \cdots \nabla_{e_{s}} m(x)^{2}$ with $s>0$ vanish if, at the end of the computation, $m^{2}$ and $\xi$ are taken constant. Again, in dimension $n=4$, the scaling dimension of the scalar field reduces to the standard $(n-2) / 2=1$.

The proof of our main Theorem 3.1 will be mainly geometric. However, we will need an intermediate analytical result, which we encapsulate in the Lemma below, which is a more detailed version of the first two paragraphs of the proof of [13, Thm. 5.1]. Logically, this analytical result follows from the Parametrized Microlocal Spectrum property, from the Locality and Covariance requirements (cf. (2) in Remark 3.2) and from the Scaling requirement. And, obviously, we make no use of either the Continuity or Analyticity requirements from [13], as we have replaced both of those by our Parametrized Microlocal Spectrum property.

Lemma 3.2. For $\left\{\tilde{\varphi}^{k}\right\}$ and $\left\{\varphi^{k}\right\}$ as in Theorem 3.1 and every fixed $M$, the identity (47) holds with some smooth functions $C_{k}[\mathbf{h}]$, where the value $C_{k}[\mathbf{h}](x)$ depends only on the germ of $\mathbf{h}$ at $x \in M$. Moreover, these functions are locally covariant, so that $\chi^{*} C_{k}[\mathbf{h}]=C_{k}\left[\chi^{*} \mathbf{h}\right]$ for any morphism $\chi$ in $\mathfrak{B k g G}$ (resp. $\mathfrak{B k g b}^{+}$), also the functions $C_{k}$ are weakly regular in the sense of Definition 2.2 and $C_{k}[\mathbf{h}]$ scales almost homogeneously of degree $k(n-2) / 2$ under the physical scaling transformation $\mathbf{h}=\left(\mathbf{g}, m^{2}, \xi\right) \mapsto\left(\lambda^{-2} \mathbf{g}, \lambda^{2} m^{2}, \xi\right)$.

Proof of Lemma 3.2. The proof is inductive in $k$. The thesis holds for $k=1$ and $C_{1}=0$, since $\varphi^{1}=\tilde{\varphi}^{1}=\varphi$. Next suppose that (47) holds for some functions $C_{i}: \Gamma(H M) \rightarrow C^{\infty}(M), i=1,2, \ldots, k-1$, that satisfy the desired properties. Then $C_{i}[\mathbf{h}] 1$ defines a locally covariant $c$-number field, where $1 \in \mathcal{W}(M, \mathbf{h})$ is the identity of the given algebra. Define

$$
\begin{equation*}
\Phi_{k,(M, \mathbf{h})}(x):=\tilde{\varphi}_{(M, \mathbf{h})}^{k}(x)-\left(\varphi_{(M, \mathbf{h})}^{k}(x)+\sum_{i=1}^{k-2}\binom{k}{i} C_{k-i}[\mathbf{h}](x) \varphi_{(M, \mathbf{h})}^{i}(x)\right) . \tag{50}
\end{equation*}
$$

By construction, $\Phi_{k}$ is a locally covariant quantum field as in Definition 3.3 and also satisfies the Algebraic, Scaling and Microlocal requirements of Definition 3.5. The algebraic properties in particular require that $\Phi_{k}$ is Hermitian and, on any given spacetime $M$, it satisfies $\left[\Phi_{k,(M, \mathbf{h})}(x), \varphi(y)\right]=0$ for all $x, y \in M$, which means that it is a $c$-number field by the kinematic completeness property of $\varphi$ (Definition 3.4). In other words, $\Phi_{k,(M, \mathbf{h})}=C_{k}[\mathbf{h}] 1$ where $C_{k}[\mathbf{h}]: C_{0}^{\infty}(M) \rightarrow \mathbb{R}$ is a distribution.

Next, we appeal to the Parametrized Microlocal Spectrum condition. That is, considering $\mathbf{h}$ itself as a 0 -parameter family, we can conclude that $C_{k}[\mathbf{h}](x)=$
$\omega\left(C_{k}[\mathbf{h}](x) 1\right)$ is a smooth function of $x$ for any Hadamard state $\omega$, since $\omega(1)=1$ for any state and Hadamard states always exist. This establishes that we have defined a map $C_{k}: \Gamma(H M) \rightarrow C^{\infty}(M)$. If we introduce an $m$-parameter family of compactly supported smooth deformation $\mathbf{H}_{s}(x)=\mathbf{H}(s, x)$ of $\mathbf{H}_{0}=\mathbf{h}$ then the same argument tells us that $C_{k}\left[\mathbf{H}_{s}\right](x)$ is also jointly smooth in $(s, x)$. Thus, according to Definition 2.2, the map $C_{k}$ is weakly regular.

The locality requirement of Definition 3.3 (see (2) in Remark 3.2) entails that $\chi^{*} C_{k}[\mathbf{h}]=C_{k}\left[\chi^{*} \mathbf{h}\right]$ for any inclusion $\chi: U \subset M$. In other words, fixing $x \in M$ and taking the limit over decreasing neighborhoods $U$ of $x$, the value $C_{k}[\mathbf{h}](x)$ depends only on the germ of $\mathbf{h}$ at $x$.

The validity of the Scaling property for both $\varphi^{k}$ and $\tilde{\varphi}^{k}$ imply that, by the formula (50), $\Phi_{k}$ is a linear combination of products of terms with almost homogeneous degrees that add up to $k(n-2) / 2$. Thus, by Lemma 2.5, $\Phi_{k}$ itself has almost homogeneous degree $k(n-2) / 2$ and thus

$$
\begin{equation*}
S_{\lambda} \Phi_{k}=\lambda^{k \frac{(n-2)}{2}} \Phi_{k}+\lambda^{k \frac{(n-2)}{2}} \sum_{i}\left(\log ^{i} \lambda\right) \Psi_{i} \tag{51}
\end{equation*}
$$

where $S_{\lambda}$ is the action of physical scalings on locally covariant scalar quantum fields, with $\Psi_{i}$ some other locally covariant quantum fields of almost homogeneous degree $k(n-2) / 2$. See Eq. (44), and the discussion below it, for the definition of $S_{\lambda}$ and in what sense locally covariant scalar quantum fields form a vector space, so that Definition 2.3 is applicable to them. Again, from the kinematic completeness of $\varphi$, it follows that $\Psi_{i,(M, \mathbf{h})}=D_{i}[\mathbf{h}] 1$ are also all $c$ number fields. On the other hand, unwrapping the definition of $S_{\lambda}$, we find that $S_{\lambda}\left(C_{k}[\mathbf{h}] 1\right)=C_{k}\left[\mathbf{h}_{\lambda}\right] 1$, and similarly for the $D_{i}$. Hence, we find that

$$
\begin{equation*}
C_{k}\left[\mathbf{h}_{\lambda}\right]=\lambda^{k \frac{(n-2)}{2}} C_{k}[\mathbf{h}]+\lambda^{k \frac{(n-2)}{2}} \sum_{i}\left(\log ^{i} \lambda\right) D_{i}[\mathbf{h}] \tag{52}
\end{equation*}
$$

is an almost homogeneous element of degree $k(n-2) / 2$ of the space of maps $\Gamma(H M) \rightarrow C^{\infty}(M)$ under the action $D \mapsto D_{\lambda}$, with $D_{\lambda}[\mathbf{h}]=D\left[\mathbf{h}_{\lambda}\right]$.

In the proof of the main Theorem below, we systematically make use of the geometric results summarized in Sect. 2. In particular, the Peetre-Slovák theorem discussed in Sect. 2.3 brings in the key simplification in our proof in comparison with the arguments of [13]. This theorem is well known in differential geometry but has not before been applied in this context. It states that, under the conditions exhibited by Lemma 3.2, the $C_{k}$ must be some (possibly non-linear) differential operators of locally bounded order applied to the background fields $\mathbf{g}, m^{2}$ and $\xi$. It then remains only to call upon the Scaling and Covariance properties to check that the $C_{k}$ may only be of the form stipulated in Eq. (48) or (49).

Proof of Theorem 3.1. In this proof, we carefully separate the hypotheses of locality, scaling and covariance. Locality allows us to conclude that the functions $C_{k}$ are differential operators. Scaling restricts their form and then covariance restricts their form even further, to the desired result. Note that, unlike in [13] we do not make use of Riemann normal coordinates. As a result, we invoke the transformations properties of $C_{k}$ under two different kinds of scaling transformations, which are mixed when normal coordinates are employed.

1. Locality and the Peetre-Slovák theorem. The first step is to combine the locality of the coefficients $C_{k}$ of Eq. (47) with the Peetre-Slovák theorem (Proposition 2.2) to conclude that in fact these coefficients are differential operators of locally bounded order (see Sect. 2.3 for details). To verify the hypotheses of Proposition 2.2, take the bundle $F \cong \mathbb{R} \times M \rightarrow M$, so that its sections are just real valued functions $\Gamma(F \rightarrow M)=C^{\infty}(M)$. Finally, take the bundle $E \cong H M \rightarrow M$. Lemma 3.2 shows that $C_{k}: \Gamma(H M) \rightarrow C^{\infty}(M)$ such that $C_{k}$ is weakly regular and $C_{k}[\mathbf{h}](x)$ depends only on the germ of $\mathbf{h}$ at $x \in M$. Consequently, the Peetre-Slovák theorem gives us the desired result: for every fixed $M \in \mathfrak{M a n}, C_{k}: \Gamma(H M) \rightarrow C^{\infty}(M)$ is a differential operator of locally bounded order, as defined in Sect. 2.3.

Although we treat $m^{2}$ and $\xi$ as spacetime-dependent fields, this is not crucial. They could be treated as constant parameters from the start and the slight modification of the proof, needed only at this point, is discussed in Remark 3.4.
2. Almost homogeneity under physical scaling. Consider a Lorentzian metric $\mathbf{g}_{0}$ on $M$, as well as a point $y \in M$ and an open neighborhood $U$ of $y$ with compact closure, with a coordinate system $\left(x^{a}\right)$ centered at $y$. Since $C_{k}$ is a differential operator of locally bounded order, for any such $\mathbf{g}_{0}, y$ and $U$ there exists ${ }^{9}$ an integer $r \geq 0$ such that $C_{k}$ is a differential operator on $U$ of local order $r$ when acting on sections of $H M$ close to ( $\mathbf{g}_{0}, m^{2}=0, \xi=0$ ), in a precise sense that we discuss next. Naturally, the coordinates $x^{a}$ induce (scaling) adapted local coordinates on the jet bundle $J^{r} H M$, which we write as $\left(x^{a}, g, g_{a b}, g^{a b, A}, w^{A}, z^{A}\right)$, recalling that the coordinates $\left(g, g_{a b}\right)$ are functionally independent up to the identity $\left|\operatorname{det} g_{a b}\right|=g$. The notation and the meaning of these coordinates are discussed in Sect. 2.2. The only difference is that we now use two sets of coordinates, $w^{A}$ and $z^{A}$, for the jets of the scalar fields, $m^{2}(x)$ and $\xi(x)$ respectively, instead of just one, and that $w$ and $z$ have corresponding scaling degrees of $s=2$ and $s=0$, as used in Sect. 2.4. Then, by the bound $r$ on the local order of $C_{k}$ at $y$, there exists a neighborhood $V_{1}^{r} \subseteq J^{r} H M$ of $j_{y}^{r}\left(\mathbf{g}_{0}, m^{2}=0, \xi=0\right)$, projecting onto $U$, and a function $F_{k}\left(x^{a}, g, g^{a b, A}, w^{A}, z^{A}\right)$ defined on $V_{1}^{r}$ such that

$$
\begin{equation*}
C_{k}[\mathbf{h}](x)=F_{k}\left(j^{r} \mathbf{h}(x)\right), \tag{53}
\end{equation*}
$$

for any section $\mathbf{h} \in \Gamma\left(\left.H M\right|_{U} \rightarrow U\right)$ such that $j^{r} \mathbf{h}(U) \subseteq V_{1}^{r}$. Note that $V_{1}^{r}$ may be strictly smaller than $\left.J^{r} H\right|_{U}$. Without loss of generality, but possibly shrinking the domain of $F_{k}$, we can choose it such that $V_{1}^{r} \cong U \times W_{1}^{r}$, where the projection on the $U$ factor is effected by the base coordinates $\left(x^{a}\right)$ and the projection onto $W_{1}^{r}$ is effected by the remaining fiber coordinates. The main obstacle to increasing $V_{1}^{r}$ to all of $J^{r} H M$ is the possible need to increase the order $r$ on larger domains. At the moment, from the Peetre-Slovák theorem, we know only that the order $r$ of $C_{k}$ is locally bounded, but may not have a finite global bound. The subscript ${ }_{1}$ on $V_{1}^{r}$ will increase in the subsequent discussion as we use the properties of $C_{k}$ to gradually enlarge the domain of definition of the function $F_{k}$, while maintaining the identity (53), and thus the bound $r$ on the order of $C_{k}$. In the final step of the proof we will in fact show

[^8]that differential order of $C_{k}$ is actually globally bounded. With that in mind, it is then consistent, on a first reading of the proof, to assume that $r$ is globally fixed and $V_{1}^{r}=J^{r} H M$, so that the parts dealing with enlarging $V^{r}$ could be skipped.

Similar to Eq. (20), the vector field implementing infinitesimal physical scaling transformations on $V_{1}^{r} \subseteq J^{r} H M$ is

$$
\begin{equation*}
e_{1}=(2+2|A|) g^{a b, A} \partial_{a b, A}+(2+2|A|) w^{A} \partial_{A}^{w}+2|A| z^{A} \partial_{A}^{z} \tag{54}
\end{equation*}
$$

According to the last statement in Lemma 3.2 and an immediate application of Lemma 2.3, the coefficient $C_{k}$ and hence the function $F_{k}$ scale almost homogeneously with degree $k(n-2) / 2$ with respect to the vector field $e_{1}$. Therefore, according to Lemma 2.4, there exists an integer $l>0$ and function $H_{j}$ on $V_{1}^{r}$, for $j=0, \ldots, l-1$, such that

$$
\begin{equation*}
F_{k}=g^{-\frac{k(n-2)}{4 n}} \sum_{j=0} \log ^{j}\left(g^{-\frac{1}{2 n}}\right) H_{j} \tag{55}
\end{equation*}
$$

where each $H_{j}$ is invariant under the action of $e_{1}$ and hence can be written as

$$
\begin{equation*}
H_{j}=H_{j}\left(x^{a}, g^{-\frac{1}{n}} g_{a b}, g^{\frac{1}{2 n}+\frac{1}{n}|A|} g^{a b, A}, g^{\frac{1}{n}+\frac{1}{n}|A|} w^{A}, g^{\frac{1}{n}|A|} z^{A}\right) \tag{56}
\end{equation*}
$$

At this point, we may extend the domain $V_{1}^{r}$ to $V_{2}^{r} \subseteq J^{r} H M$, which is invariant under physical scaling. That is, we can write $V_{2}^{r} \cong \mathbb{R}^{+} \times W_{2}^{r}$, where the coordinate $g$ effects the projection onto the $\mathbb{R}^{+}$factor and the coordinates $\left(x^{a}, g^{-\frac{1}{n}} g_{a b}, g^{\frac{1}{2 n}+\frac{1}{n}|A|} g^{a b, A}, g^{\frac{1}{n}+\frac{1}{n}|A|} w^{A}, g^{\frac{1}{n}|A|} z^{A}\right)$ effect the projection onto the $W_{2}^{r}$ factor, which includes at least the point $\left(g^{-\frac{1}{n}} g_{a b} \circ \mathbf{g}_{0}(y), 0,0,0\right)$. The function $F_{k}$ extends from $V_{1}^{r}$ to $V_{2}^{r}$ in a unique way as an almost homogeneous function of degree $k(n-2) / 2$.

Let us go into some of the details of the mentioned unique extension procedure. So far, we could only presume that the identity (53) that expresses the function $C_{k}[\mathbf{h}](x)$ in terms of the differential operator defined by the function $F_{k}$ holds only when the germ of $\mathbf{h}$ at $x \in M$ projects onto one of the jets in the domain $V_{1}^{r} \subseteq J^{r} H M$ of $F_{k}$. We have defined the extended domain $V_{2}^{r}$ to be the smallest domain invariant under physical scaling and containing $V_{1}^{r}$. The function $F_{k}$, by using formula (55), has a unique almost homogeneous extension to $V_{2}^{r}$ that scales almost homogeneously and agrees with the known values of $F_{k}$ on $V_{1}^{r}$. Since any element of $V_{2}^{r}$ can be brought back to $V_{1}^{r}$ by a physical scaling transformation and $C_{k}[\mathbf{h}]$ itself scales almost homogeneously, the identity (53) must remain valid also for germs of $\mathbf{h}$ at $x$ that project to jets in the extended domain $V_{2}^{r}$. Below, we use similar logic each time the domain of the function $F_{k}$ is expanded, eventually to all of $J^{r} H M$, though possibly with a larger value of $r$, thus showing that $C_{k}[\mathbf{h}]$ is actually a differential operator of globally bounded order.
3. Diffeomorphism covariance and the Thomas replacement theorem. Now we move on to the covariance property of the $C_{k}$ under diffeomorphisms, which will be used in several stages. First, fixing the previously made choice of $y \in M$, we note that the preceding arguments using the Peetre-Slovák theorem can be repeated for any pair of $y^{\prime} \in M$ and $\mathbf{g}_{0}^{\prime}=\chi^{*} \mathbf{g}_{0}$, where $\chi: M \rightarrow M$ is some diffeomorphism such that $\chi\left(y^{\prime}\right)=y$, giving rise to differential orders $r^{\prime}$ and domains $V_{2}^{\prime r^{\prime}} \subseteq J^{r^{\prime}} H M$. The diffeomorphism covariance of $C_{k}$ then implies
that all these differential orders are the same, $r^{\prime}=r$, and that the union $V_{3}^{\prime r} \subseteq$ $J^{r} H M$ of all the $V_{2}^{\prime r^{\prime}}$ domains defines a neighborhood of the $\operatorname{Diff}(M)$-orbit of $j^{r}\left(\mathbf{g}_{0}, 0,0\right) \in J^{r} H M$. In fact, $V_{3}^{\prime r}$ can itself be chosen to be Diff( $M$ )-invariant (for instance, by taking the union of all $\operatorname{Diff}(M)$ images of a non-invariant $V_{3}^{\prime r}$ ) and a function $F_{k}$ satisfying (53) uniquely defined on it. The diffeomorphism covariance of $C_{k}$ then implies that $F_{k}$ is itself $\operatorname{Diff}(M)$-invariant on $V_{3}^{\prime r}$, in the sense described in Sect. 2.5. The case of $\operatorname{Diff}^{+}(M)$ covariance is handled in exactly the same way.

Since diffeomorphisms act transitively on $M$, a diffeomorphism invariant $V_{3}^{\prime r}$ would then project down to all of $M$. Instead, motivated by the desire to keep working in the coordinates adapted to the local chart $\left(x^{a}\right)$ on $U \subseteq M$, we choose $V_{3}^{r}$ instead to be the intersection of $V_{3}^{\prime r}$ and the pre-image of $\bar{U}$ under the projection $J^{r} H M \rightarrow M$. Then, we have all the needed hypothesis to apply Proposition 2.6 to eliminate the dependence of $F_{k}$, as a diffeomorphism invariant function, on some of the coordinates on $V_{3}^{r}$. Actually, part of the almost homogeneous scaling property implies that the functions $H_{j}$ from Eq. (56) are each separately invariant under diffeomorphisms, so that we can apply Proposition 2.6 to each of them individually. Therefore, we can conclude that

$$
\begin{align*}
& g^{-\frac{k(n-2)}{4 n}} H_{j}\left(x^{a}, g^{-\frac{1}{n}} g_{a b}, g^{\frac{1}{n}+\frac{1}{n}|A|} g^{a b, A}, g^{\frac{1}{n}+\frac{1}{n}|A|} w^{A}, g^{\frac{1}{n}|A|} z^{A}\right) \\
& \quad=g^{-\frac{k(n-2)}{4 n}} G_{j}\left(g^{-\frac{1}{n}} g_{a b}, g^{\frac{3}{n}+\frac{1}{n}|A|} \bar{S}^{a b(c d, A)}, g^{\frac{1}{n}+\frac{1}{n}|A|} \bar{w}^{A}, g^{\frac{1}{n}|A|} \bar{z}^{A}\right) \tag{57}
\end{align*}
$$

where the notation used for the coordinates is explained in Sect. 2.5 and each $g^{-\frac{k(n-2)}{4 n}} G_{j}$, for $j=0, \ldots, l-1$, is invariant under the natural action of either $G L(n)$ (or $G L^{+}(n)$, depending on which of the cases (a) or (b) we are dealing with) on its arguments. Notably, $G_{j}$ depends neither on the base ( $x^{a}$ ) nor on the Christoffel coordinates $\left(\Gamma_{(b c, A)}^{a}\right)$.

The invariance properties of $V_{3}^{r}$ now tells us that it has the structure $V_{3}^{r} \cong$ $U \times L_{n} \times \mathbb{R}^{\gamma} \times W_{3}^{r}$, where the coordinates $\left(x^{a}\right)$ effect the projection onto the $U$ factor, the coordinates $\left(g_{a b}\right)$ or $\left(g, g^{-\frac{1}{n}} g_{a b}\right)$ effect the projection onto the $L_{n}$ factor (the whole space of non-degenerate bilinear forms on $\mathbb{R}^{n}$ with Lorentzian signature), the coordinates $\Gamma_{(b c, A)}^{a}$ effect the projection on the $\mathbb{R}^{\gamma}$ argument and the remaining coordinates $\left(g^{\overline{3} n+\frac{1}{n}|A|} \bar{S}^{a b(c d, A)}, g^{\frac{1}{n}+\frac{1}{n}|A|} \bar{w}^{A}, g^{\frac{1}{n}+\frac{1}{n}|A|} \bar{z}^{A}\right)$ effect the projection on the $W_{3}^{r}$ argument, which contains at least the point $(0,0,0)$ and is invariant under the corresponding action of $G L(n)$ (resp. $\left.G L^{+}(n)\right)$.
4. Invariance under coordinate scaling. Next, recall the action of the subgroup of $G L(n)\left(\right.$ resp. $\left.G L^{+}(n)\right)$ that we called coordinate scalings in Sect. 2.5. Notice that all the coordinates that the functions $G_{j}$ depend on have positive weight with respect to coordinate scalings, with the exception of $\left(g^{-\frac{1}{n}} g_{a b}, z\right)$. For brevity, let us rewrite our coordinates as $\left(g, g^{-\frac{1}{n}} g_{a b}, z, q^{i}\right)$, with the weight of the coordinate $q^{i}$ under coordinate scalings denoted by $d_{i}>0$. Then the invariance of the functions $F_{k}$ on $V_{3}^{r}$ under diffeomorphisms, and hence coordinate scalings, implies the identity

$$
\begin{align*}
\mu^{k \frac{(n-2)}{2}} F_{k}\left(g, g^{-\frac{1}{n}} g_{a b}, z, q^{i}\right) & =\mu^{k \frac{(n-2)}{2}} F_{k}\left(\mu^{2 n} g, g^{-\frac{1}{n}} g_{a b}, z, \mu^{d_{i}} q^{i}\right) \\
& =g^{-\frac{k(n-2)}{4 n}} \sum_{j=0}^{l-1} \log ^{j}\left(\mu^{-1} g^{-\frac{1}{2 n}}\right) G_{j}\left(g^{-\frac{1}{n}} g_{a b}, z, \mu^{d_{i}} q^{i}\right) \tag{58}
\end{align*}
$$

for any point of $V_{3}^{r}$ on its left hand side and any value of $\mu>0$. As described above, the limit $\left(g^{-\frac{1}{n}} g_{a b}, z, 0\right)$ of the arguments of the functions $G_{j}$ as $\mu \rightarrow 0$ falls within the domain of the functions $G_{j}$. Therefore, while the limit of the left-hand side of (58) converges to 0 as $\mu \rightarrow 0$, the right-hand side diverges unless all $G_{j}=0$ for $j>0$, so that $F_{k}=g^{-\frac{k(n-2)}{4 n}} G_{0}$. The new identity implied by invariance under coordinate scalings is then

$$
\begin{equation*}
g^{-\frac{k(n-2)}{4 n}} G_{0}\left(g^{-\frac{1}{n}} g_{a b}, z, q^{i}\right)=\mu^{-k \frac{(n-2)}{2}} g^{-\frac{k(n-2)}{4 n}} G_{0}\left(g^{-\frac{1}{n}} g_{a b}, z, \mu^{d_{i}} q^{i}\right) \tag{59}
\end{equation*}
$$

Fix some values for the coordinates $\left(g, g^{-\frac{1}{n}} g_{a b}, z\right)$ and recall that the point $\left(g^{-\frac{1}{n}} g_{a b}, z, 0\right)$ is part of the domain of definition of $G_{0}$. Since $G_{0}$ is smooth, Taylor's theorem allows us to write it as

$$
\begin{equation*}
G_{0}\left(g^{-\frac{1}{n}} g_{a b}, z, q^{i}\right)=\sum_{|I|<N} A_{I}\left(g^{-\frac{1}{n}} g_{a b}, z\right) q^{I}+O\left(q^{N}\right) \tag{60}
\end{equation*}
$$

where $I=i_{1} \cdots i_{m}$ is a multi-index with respect to the coordinates ( $q^{i}$ ) and $N>0$ is an integer large enough so that $\langle d, I\rangle=\sum_{j=1}^{m} d_{i_{j}}>k$ for any $m=$ $|I|>N$. Note that the error term $O\left(q^{N}\right)$, for fixed $\left(q^{i}\right)$ mapped to $\left(\mu^{d_{i}} q^{i}\right)$ and $\mu \rightarrow 0$, is mapped to $O\left(\mu^{k+1}\right)$ by our choice of sufficiently large $N$. Thus, using Taylor's theorem, we can rewrite (59) as

$$
\begin{align*}
g^{-\frac{k(n-2)}{4 n}} G_{0}\left(g^{-\frac{1}{n}} g_{a b}, z, q^{i}\right)=\sum_{|I|<N} g^{-\frac{k(n-2)}{4 n}} A_{I}( & \left.g^{-\frac{1}{n}} g_{a b}, z\right) q^{I} \mu^{\langle d, I\rangle-k \frac{(n-2)}{2}} \\
& +\mu^{-\frac{k(n-2)}{2}} O\left(\mu^{\frac{k(n-2)}{2}+1}\right) \tag{61}
\end{align*}
$$

While the left-hand side of (61) is bounded as $\mu \rightarrow 0$, the right-hand side diverges unless all $A_{I}=0$ for $I$ such that $\langle d, I\rangle<k(n-2) / 2$. If this vanishing condition is satisfied, the $\mu \rightarrow 0$ limits of both sides of (61) exist and give the identity

$$
\begin{equation*}
F_{k}=g^{-\frac{k(n-2)}{4 n}} G_{0}\left(g^{-\frac{1}{n}} g_{a b}, z, q^{i}\right)=\sum_{\langle d, I\rangle=\frac{k(n-2)}{2}} g^{-\frac{k(n-2)}{4 n}} A_{I}\left(g^{-\frac{1}{n}} g_{a b}, z\right) q^{I} \tag{62}
\end{equation*}
$$

At this point, we can once more enlarge the domain of definition of the function $F_{k}$, where the identity (53) holds, from $V_{3}^{r}$ to $V_{4}^{r} \subset J^{r} H M$. The new domain is isomorphic to $V_{4}^{r} \cong U \times L_{n} \times \mathbb{R} \times W_{4} \times \mathbb{R}^{\gamma} \times \mathbb{R}^{\delta}$, where the coordinates $\left(x^{a}\right)$ effect the projection onto the $U$ factor, the coordinates $\left(g_{a b}\right)$ or $\left(g, g^{-\frac{1}{n}} g_{a b}\right)$ effect the projection onto the $L_{n}$ factor, the coordinate $\left(g^{\frac{1}{n}} w\right)$ effects the projection onto the $\mathbb{R}$ factor, the coordinate $(z)$ effects the projection onto the $W_{4}$ factor (which at least contains the point (0)), the coordinates $\left(\Gamma_{(b c, A)}^{a}\right)$ effect the projection onto the $\mathbb{R}^{\gamma}$ factor, and the remaining coordinates $\left(g^{\frac{3}{n}+\frac{1}{n}|A|} \bar{S}^{a b(c d, A)}, g^{\frac{1}{n}+\frac{1}{n}|A|} \bar{w}^{A}, g^{\frac{1}{n}+\frac{1}{n}|A|} \bar{z}^{A}\right)$ effect the projection onto the $\mathbb{R}^{\delta}$ factor, where the coordinates involving $\bar{w}^{A}$ and $\bar{z}^{A}$ with $|A|=0$ are obviously excluded. Note that $U \times L_{n} \times \mathbb{R} \times W_{4} \subseteq H M$ and that $V_{4}^{r}$ is simply its preimage with respect to the bundle projection $J^{r} H M \rightarrow H M$. The function $F_{k}$ extends uniquely from $V_{3}^{r}$ to a function on $V_{4}^{r}$ that is invariant under coordinate scalings. The reason we could extend the domain so much, essentially the factor $W_{3}^{r}$ got enlarged to $\mathbb{R} \times W_{4} \times \mathbb{R}^{\delta}$, is because almost all coordinates, those we
labeled by $\left(q^{i}\right)$ above, had positive degrees with respect to coordinate scalings. The range of the $(z)$ coordinate is limited to $W_{4}$ because it is invariant under coordinate scalings and even under the larger group $G L(n)$ (resp. $G L^{+}(n)$ ) that acts on the other bundle coordinates. Also, note that according to Eq. (62) the dependence of $F_{k}$ on the $\mathbb{R} \times \mathbb{R}^{\gamma} \times \mathbb{R}^{\delta}$ factor in $V_{4}^{r}$, corresponding to the coordinates we labeled by ( $q^{i}$ ) above, is polynomial.
5. $G L(n)$-equivariance and polynomial dependence on the metric. From the preceding discussion, the function $F_{k}$, satisfying the identity (53), is defined on the domain $V_{4}^{r}=U \times V_{4}$ and depends only on the coordinates corresponding to the factor $V_{4}=L_{n} \times W_{4} \times \mathbb{R}^{\delta}$ (where we have grouped all the $\mathbb{R} \times \mathbb{R}^{\gamma} \times \mathbb{R}^{\delta}$ factors together into $\mathbb{R}^{\delta}$, implicitly redefining $\delta$ ). Moreover, the dependence on the coordinates on the $\mathbb{R}^{\delta}$ factor is polynomial, while the coefficients $g^{-\frac{k(n-2)}{4 n}} A_{I}\left(g^{-\frac{1}{n}} g_{a b}, z\right)$ of these polynomials depend only on the $L_{n} \times W_{4}$ factor. It is also clear from the preceding discussion that each of the factors in $V_{4}$ carries a tensor density representation of $G L(n)\left(\right.$ resp. $\left.G L^{+}(n)\right)($ cf. Sect. 2.6), which happens to be trivial on $W_{4}$. The space of functions on $V_{4}$ then itself carries a representation of $G L(n)$ (resp. $G L^{+}(n)$ ), induced by the pullback of the action on $V_{4}$, and the function $F_{k}$ is invariant under this action. In the same way, the space $\mathcal{P}_{\delta}^{N}$ of polynomials of degree no greater than $N$ on $\mathbb{R}^{\delta}$ carries a representation of $G L(n)$ (resp. $G L^{+}(n)$ ),

$$
\begin{equation*}
(u P)(\rho)=P\left(u^{-1} \rho\right), \quad \text { for any } u \in G L(n), P \in \mathcal{P}_{\delta}^{N} \text { and } \rho \in \mathbb{R}^{\delta} \tag{63}
\end{equation*}
$$

which by elementary reasoning, within the representation theory of $G L(n)$ [9], is a direct sum of tensor density representations. Let us group these subrepresentations by tensor rank and density weight. Therefore, $\mathcal{P}_{\delta}^{N}=\bigoplus_{j} T_{j}$, where each $T_{j}$ is a tensor density representation.

The form that we have reduced $F_{k}$ to can be described as follows. Given a point $(\mathbf{g}, \xi, \rho) \in V_{4}$, the $A$-coefficients $g^{-\frac{k(n-2)}{4 n}} A_{I}\left(g^{-\frac{1}{n}} g_{a b}, z\right)$ evaluated at $(\mathbf{g}, \xi) \in L_{n} \times W_{4}$ give a polynomial in $\mathcal{P}_{\delta}^{N}$, which is then evaluated at $\rho \in \mathbb{R}^{\delta}$. Thus we can think of the $A$-coefficients as a collection of functions $A_{j}: L_{n} \times$ $W_{4} \rightarrow T_{j}$, with components given by

$$
\begin{equation*}
\left(A_{j}\left(g_{a b}, z\right)\right)_{I}=g^{-\frac{k(n-2)}{4 n}} A_{I}\left(g^{-\frac{1}{n}} g_{a b}, z\right) \tag{64}
\end{equation*}
$$

The only way for $F_{k}$ constructed in this way to be invariant under the action of $G L(n)$ is for the maps $A_{j}$ to be equivariant (cf. Sect. 2.6), so that

$$
\begin{align*}
F_{k}(u \mathbf{g}, u \xi, u \rho)= & \sum_{j} A_{j}(u \mathbf{g}, u \xi)(u \rho)=\sum_{j}\left(u A_{j}(\mathbf{g}, \xi)\right)(u \rho) \\
& =\sum_{j} A_{j}(\mathbf{g}, \xi)\left(u^{-1} u \rho\right)=\sum_{j} A_{j}(\mathbf{g}, \xi)(\rho)=F_{k}(\mathbf{g}, \xi, \rho) \tag{65}
\end{align*}
$$

for any $u \in G L(n)\left(\right.$ resp. $\left.G L^{+}(n)\right)$ and $(\mathbf{g}, \xi, \rho) \in V_{4}$.
We are finally in a position to conclude that, for a fixed $\xi \in W_{4}$, the map $A_{j}(-, \xi): L_{n} \rightarrow T_{j}$ is an equivariant tensor density, in the sense of Definition 2.7, and hence must be of the form dictated by Lemma 2.8, which characterizes all such maps in a way, in view of Remark 2.5, compatible with our formula (64). In other words, the coefficients of the polynomials $A_{j}(\mathbf{g}, \xi)$ depend themselves polynomially on the components $g_{a b}$ and $\varepsilon_{a_{1} \cdots a_{n}}$ of the covariant metric and

Levi-Civita tensors, up to an overall multiple of $g=\left|\operatorname{det} g_{a b}\right|$. If $F_{k}$ is invariant under $G L(n)$, then the dependence on $\varepsilon_{a_{1} \cdots a_{n}}$ must be trivial, while it could in general be non-trivial if $F_{k}$ is invariant only under $G L^{+}(n)$. Expanding all the polynomials in $g_{a b}, \varepsilon_{a_{1} \cdots a_{n}}$ and $q^{i}$, all the factors of powers of $g$ must collectively cancel to preserve invariance of $F_{k}$ under $G L(n)$ (resp. $G L^{+}(n)$ ). In other words, we can conclude that

$$
\begin{equation*}
F_{k}=\sum_{j} a_{j}(z) P_{j}\left(g_{a b}, \varepsilon_{a_{1} \cdots a_{n}}, \bar{S}^{a b(c d, A)}, \bar{w}^{A}, \bar{z}^{A}\right), \quad \text { with }|A| \geq 1 \text { in } \bar{z}^{A} \tag{66}
\end{equation*}
$$

where the sum is over a (necessarily finite) basis of polynomials $P_{j}$, which consist of linear combinations of tensor contractions of products of their arguments, with coefficients arbitrarily depending on the $z$ coordinate. In this form, the function $F_{k}$ is manifestly invariant under $G L(n)$ (resp. $G L^{+}(n)$ ) transformations.
6. Global boundedness of differential order. To conclude the proof, it remains only to extend the domain $V_{4}^{r}$ once more, this time to all of $J^{r} H M$, for an appropriate choice of $r$. It is well known that for a fixed weight $s$ under physical scaling, there is only a finite number of linearly independent polynomials $P_{j}$ of weight $s$ constructed, as described above, from the metric and the covariant derivatives of the scalar fields $m^{2}, \xi$ and the Riemann curvature tensor in the form $\bar{S}^{a b c d}$, even if the number of the derivatives $r$ is allowed to be arbitrary ${ }^{10}$ [8]. Let $r_{k}$ be the maximum number of derivatives that appear in a basis for these polynomials $P_{j}$ when $s=k(n-2) / 2$. Then, no matter the original choice of domain $U \subseteq M$, the differential operator $C_{k}$ restricted to it must be of order $\leq r_{k}$. Thus, we are justified in setting $r=r_{k}$ in all of the preceding discussion. The only obstacle that may have prevented us from extending the domain $V_{4}^{r} \subseteq J^{r_{k}} H M$ of the function $F_{k}$ to all of the pre-image of $U$ under the projection $J^{r_{k}} H M \rightarrow M$ is the possibility that $C_{k}$ would change order on jets whose projections fall outside $V_{4}^{r_{k}}$. However, with the maximal possible order of $C_{k}$ bounded by $r_{k}$, this obstacle is now absent. In other words, we can safely presume that $V_{4}^{r_{k}}$ is equal to the pre-image of $U \subseteq M$ with respect to the projection $J^{r_{k}} H M \rightarrow M$, with $F_{k}$ retaining the form (66) on all of its domain. A slightly more detailed version of this argument would note that the original choice of the domain $V_{1}^{r}$ to be a neighborhood of a point $j_{y}^{r}\left(\mathbf{g}_{0}, m^{2}=0, \xi=0\right)$ in $J^{r} H M$ could have equally been chosen to be a neighborhood of the point $j_{y}^{r}\left(\mathbf{g}_{0}, m^{2}=0, \xi_{0}\right)$, without affecting any subsequent arguments. Piecing together $F_{k}$ over the extensions of all such neighborhoods gives us a definition of $F_{k}$ on the entire pre-image of $U$ under the projection $J^{r_{k}} H M \rightarrow M$ with the same global order bound $r_{k}$. Finally, covariance of $C_{k}$ with respect to diffeomorphisms requires that the form (66) is also independent of the domain $U \subseteq M$. Thus, we can conclude that there exists a globally defined smooth bundle map

$$
\begin{aligned}
& { }^{10} \text { To see that, consider a monomial of the schematic form } \\
& \qquad\left(g_{a b}\right)^{p_{g}}\left(\varepsilon_{a_{1} \cdots a_{n}}\right)^{p_{\varepsilon}} \prod_{|A|}\left(\bar{S}^{a b(c d, A)}\right)^{p_{S,|A|}}\left(\bar{w}^{A}\right)^{p_{w,|A|}}\left(\bar{z}^{A}\right)^{p_{z,|A|}}
\end{aligned}
$$

necessarily with $p_{z, 0}=0$, and note that the $p$-exponents must satisfy the constraint $\sum_{|A|}[(2+$ $\left.|A|) p_{S,|A|}+(1+|A|) p_{w,|A|}+|A| p_{z,|A|}\right]=s$, due to $s$-homogeneity with respect to physical scalings and invariance with respect to coordinate scalings. Since each $p$-exponent is nonnegative, this implies a bound on the maximum value of $|A|$ with a non-zero exponent.
$F_{k}: J^{r_{k}} H M \rightarrow \mathbb{R} \times M$ over $M$ of the form (66) such that $C_{k}[\mathbf{h}](x)=F_{k} \circ j^{r_{k}} \mathbf{h}(x)$ for any $x \in M$ and $\mathbf{h} \in \Gamma(H M)$, which concludes the proof.

Remark 3.4. We observe, by looking at the Locality and the Peetre-Slovák theorem step of the above proof and also at the proof of Lemma 3.2, that one might wonder at the need to take $m^{2}$ and $\xi$ as spacetime-dependent fields rather than constants, as is usually the case. Our arguments still go through, with only two changes. First, the microlocal hypothesis mentioned in 3.5 must be strengthened to require an empty wavefront set for $\omega\left(\varphi^{k}(x)\right)$ as a distribution on $M \times \mathbb{R}^{2}$ (with the $\mathbb{R}^{2}$ factor standing for the parameter space of $m^{2}$ and $\xi$ ) rather than as a distribution on $M$ for any fixed $m^{2}$ and $\xi$. Note that the weaker microlocal requirement does not exclude the infinite family of counterterms of [25] that were discussed in the Introduction, while the stronger one does. Second, we must make use of the more general version of the Peetre-Slovák theorem for differential operators with parameters, as in Proposition A. 1 in Appendix A. To apply that result, we would need to let $N=M$ and replace the spacetime manifold $M$ by $P=M \times \mathbb{R}^{2}$, adding the $\left(m^{2}, \xi\right)$ parameter space. It would then follow from known information about $C_{k}$ that it is local with respect to the natural projection $P \cong \mathbb{R}^{2} \times M \rightarrow M$, hence satisfying the more general Peetre-Slovák theorem.

We end this section with a couple of straight forward but noteworthy observations. First, it is a direct result of the proof of Lemma 3.2 that the set of coefficients $\left\{C_{k}[\mathbf{h}]\right\}$ from Eq. (47) is determined jointly by the entire families $\left\{\varphi^{k}\right\}$ and $\left\{\tilde{\varphi}^{k}\right\}$ of Wick powers, rather than depending on each pair $\varphi^{k}$ and $\tilde{\varphi}^{k}$ individually. Second, the converse of Theorem 3.1 holds as well. That is, given a family $\left\{\varphi^{k}\right\}$ of locally covariant Wick powers and a set $\left\{C_{k}[\mathbf{h}]\right\}$ of satisfying the conclusions of Theorem 3.1, the formula (47) defines another family $\left\{\varphi^{k}\right\}$ of locally covariant Wick powers.

## 4 Discussion

In this work, we have characterized admissible finite renormalizations of Wick powers of a locally covariant quantum scalar field $\varphi$ on curved spacetimes, with possibly spacetime-dependent mass $m^{2}$ and curvature coupling $\xi$. By local covariance, we mean the axioms of Brunetti, Fredenhagen and Verch [7]. Our work is a significant technical improvement on the original work of Hollands and Wald [13] on this subject. The main result (Theorem 3.1) is a slight generalization of that of Hollands and Wald, yet our hypotheses are significantly more natural and the proof is greatly simplified and streamlined.

Under standard hypotheses, on Minkowski space, where the curvature coupling $\xi$ is absent, it is well known that the finite renormalizations of the Wick powers $\varphi^{k}$ are restricted to linear combinations of Wick powers of lower order, with dimensionful coefficients that are polynomials in $m^{2}$, with the total dimension matching that of $\varphi^{k}$. This is a strong constraint, because the resulting space of possibilities is finite-dimensional. On curved spacetimes, as first proven by Hollands and Wald in [13], adding local covariance and some further more technical hypotheses gives a result of comparable strength. The only modification is that the coefficients of lower order Wick powers can also depend polynomially
on curvature scalars and analytically on $\xi$, with the same restriction on their dimensions. The resulting possibilities no longer form a finite-dimensional space, but a quasi-finite-dimensional one, in the sense that it is finitely generated under linear combinations with coefficients analytic in $\xi$. It is worth noting that the dependence of finite renormalization terms on the background metric is entirely contained in the curvature scalars, while their $\xi$-dependent coefficients must be assigned uniformly across all spacetimes to preserve local covariance.

The hypotheses of Hollands and Wald, briefly recalled in Definition 3.5, include the requirements of locality, microlocal regularity and of continuous and analytic dependence on the background spacetime metric and coupling parameters. Unfortunately, while playing a crucial role in the existing proof, the analytic dependence hypothesis has been long considered somewhat unnatural and technically very cumbersome. We have found that, by using a standard result of differential geometry (the Peetre-Slovák theorem, cf. Proposition 2.2 and Appendix A), in the presence of the remaining assumptions, the role of both the continuity and analyticity hypotheses is completely subsumed by that of locality and a strengthened version of the microlocal regularity condition. We believe the strengthened, so-called microlocal spectral condition is natural from both physical and geometrical points of view. Physically it encapsulates the stability of the microlocal properties of Wick powers under smooth variations of the background geometry. Geometrically, it provides just the right hypothesis needed to prove the locality of finite renormalizations of Wick powers, without leaving the realm of smooth differential geometry. Thus, by replacing the continuity and analyticity requirements by a more natural hypothesis, our final result on the characterization of finite renormalizations of Wick powers, as stated in Theorem 3.1, is essentially identical to that of Hollands and Wald. The main difference is that arbitrary smooth dependence on the coupling $\xi$ is now allowed, instead of just analytic dependence. Another difference is that we have also explicitly considered weakening covariance to only under orientation preserving diffeomorphisms, which increases the renormalization freedom to curvature scalars constructed also with the Levi-Civita tensor and not just the metric. Finally, we explicitly treat $m^{2}$ and $\xi$ as possibly spacetime-dependent parameters, rather than simple constants. The original proof of Hollands and Wald also treated them as spacetime-dependent, while restricting to the case of constants in the statement of their final result. We noted in Remark 3.4 how our arguments could be adapted to treating the parameters as constants throughout.

As was already mentioned, our characterization of finite renormalizations extends to theories that need only be covariant with respect to orientation preserving diffeomorphisms. In particular, in even dimensions, chiral theories (those not invariant under spatial parity transformations) could be admissible. While our result does not contain any surprises, it is important to have a rigorous statement on the complete range of possibilities. In particular, suppose that a classical parity invariant theory is perturbatively quantized using a chiral renormalization scheme. The knowledge of a complete classification of finite renormalizations is then required to decide whether there exists a different renormalization scheme that gives a parity invariant quantization.

Another advantage of our proof is the clear separation between the applications of the locality, microlocal regularity, covariance and scaling hypotheses. We make a particular distinction between physical scalings (those resulting from
a rescaling of the metric) and coordinate scalings (those resulting from the local action of some diffeomorphisms). We believe that structuring the proof in this way makes it significantly easier to generalize the result to other types of tensor or spinor fields, a task that is yet to be seriously taken up in the literature on locally covariant quantum field theory, which is in significant part likely due to the complexity of and the unnatural hypotheses needed in the original proof of Hollands and Wald. In particular, it is likely that the crucial step in limiting the finite renormalization freedom to a quasi-finite-dimensional space is to carefully balance the covariance and scaling properties, such that there exists a coordinate system on the jets of background fields, like the rescaled curvature coordinates that we identified in Sect. 2.2, where all coordinates corresponding higher derivatives have positive weight under a combination of the physical and coordinate scalings.

Another direction in which our main result could be generalized is to consider Wick powers that included derivatives of fields. Our proof should extend without problems. The main difference would be that the finite renormalization coefficients $C_{k}$ could then be tensor- instead of scalar-valued, since Wick powers with derivatives could themselves be tensor fields. This difference would affect the part of our proof where we make use of $G L(n)$-equivariance to fix the form of the $A_{j}$ coefficients, which could be mixed densitized tensors. Fortunately, the main technical result on the classification of equivariant tensor densities, as stated in Lemma 2.8, is sufficiently general to apply to that case as well, since introducing densitization erases the distinction between covariant and contravariant indices.

Let us also say something about time-ordered products. Hollands and Wald also gave a sketch of the proof of the characterization of finite renormalizations of time-ordered products [13, Thm. 5.2], under the same hypothesis as their result about Wick powers. As they point out, the main difference with the case of Wick powers is in the structure of coefficients that are analogous to the $C_{k}$, which become distributions on multiple copies of the spacetime manifold. The arguments, which we encapsulated in Lemma 3.2, applying microlocal arguments to restrict the wavefront set of these distributions would have to be generalized accordingly. After that point, the proof of Theorem 3.1, would apply without essential modifications. Thus, our methods should generalize to time-ordered products as well.

Finally, we note that, although we believe that our parametrized microlocal spectral condition (Definition 3.5(iv)) does hold for standard locally covariant Hadamard parametrix prescription for defining Wick powers, we have not given a proof. In fact a complete proof of the validity of the continuous and analytic dependence for the Hadamard parametrix prescription was not given in the original work of Hollands and Wald either [13, Sec. 5.2] and only appeared in the later work [14]. Closing this gap with a complete and precise proof is a worthwhile goal for future work. In fact, to be of greatest use for the characterization of finite renormalizations of time-ordered products and of Wick products with derivatives, we would need a proof of validity of a parametrized microlocal spectral condition for multi-local fields as well. Such a condition might be reasonably stated as follows, echoing [13, Eq. (46)] which considered a similar question for the analytic wavefront set. Let $\Phi\left(x_{1}, \ldots, x_{k}\right)$ be a locally covariant multi-local field that already satisfies the standard (unparametrized) microlocal spectral condition. Then, let $(M, \mathbf{h})$ be a background geometry and
let

$$
\begin{equation*}
\Gamma^{\Phi}\left(M^{k}, \mathbf{h}\right)=\overline{\bigcup_{\omega} \mathrm{WF}\left(\omega\left(\Phi_{(M, \mathbf{h})}(-)\right) \backslash\{0\} \subseteq\left(T^{*} M^{k}\right) \backslash\{0\}, \text {, },\right. \text {. }} \tag{67}
\end{equation*}
$$

where the union is taken over all Hadamard states $\omega$ on $\mathcal{W}(M, \mathbf{h})$. Consider also a smooth $m$-parameter compactly supported variation $\mathbf{H}(s, x)$ of $\mathbf{h}(x)$, together with the accompanying algebra isomorphisms $\tau_{\text {ret }}^{s}: \mathcal{W}\left(M, \mathbf{H}_{s}\right) \rightarrow \mathcal{W}(M, \mathbf{h})$. If $\omega$ is any Hadamard state on $\mathcal{W}(M, \mathbf{h})$, then we would like to require that the wavefront set of $E_{\omega}^{\Phi}\left(s, x_{1}, \ldots, x_{k}\right)=\omega\left(\tau_{\text {ret }}^{s} \circ \Phi_{\left(M, \mathbf{H}_{s}\right)}\left(x_{1}, \ldots, x_{k}\right)\right)$ as a distribution on $\mathbb{R}^{m} \times M^{k}$ satisfies

$$
\begin{align*}
\mathrm{WF}\left(E_{\omega}^{\Phi}\right) \subseteq\left\{\left(s, \sigma, x_{1}, p_{1}, \ldots, x_{n}, p_{n}\right)\right. & \in T^{*}\left(\mathbb{R}^{m} \times M^{k}\right) \\
& \left.\mid\left(x_{1}, p_{1}, \ldots, x_{k}, p_{k}\right) \in \Gamma^{\Phi}\left(M, \mathbf{H}_{s}\right)\right\} . \tag{68}
\end{align*}
$$

That is, $E_{\omega}^{\varphi}$ or any of its derivatives can be restricted to the submanifold of $\mathbb{R}^{m} \times M^{k}$ given by a fixed value of $s$, and that restriction has precisely the wavefront set expected of a locally covariant field on $\left(M, \mathbf{H}_{s}\right)$ satisfying the standard microlocal spectrum condition.

We now give a sketch of an argument for establishing the validity of our parametrized microlocal spectrum condition in the simplest case of the Wick monomial $\phi_{H}^{2}(x)$, as constructed using the Hadamard parametrix regularization method by Hollands and Wald. (The authors are grateful to an anonymous referee for suggesting it.) In this elementary situation, the only thing to prove is the joint smoothness of $\omega_{s}\left(\phi_{H_{s}}^{2}(x)\right)$ as a function of the position $x$ and the parameters $s$, with $H_{s}$ the local Hadamard parametrices of the compactly supported background geometry variation $\mathbf{H}_{s}$ and the Hadamard states $\omega_{s}=\omega \circ \tau_{\text {ret }}^{s}$ are defined as in the preceding paragraph. As a matter of fact, one could more generally prove that $f(s, x, y):=\omega_{s}(\phi(x) \phi(y))-H_{s}(x, y)$ is a jointly smooth function of $s, x, y$, since $\omega_{s}\left(\phi_{H_{s}}^{2}(x)\right)$ is obtained just by taking $x=y$ in the difference above. The 2-point function $\omega_{s}(\phi(x) \phi(y))$ is a global distributional bisolution, while the symmetric distribution $H_{s}(x, y)$ is locally defined through a well-known procedure and, as a parametrix, only satisfies the Klein-Gordon equation (in either $x$ or $y$ ) up to a smooth error term, which we would need to show is also jointly smooth in $s$. So, the function $f(s, x, y)$ will satisfy a pair of Klein-Gordon equations,

$$
\begin{align*}
\left(\square_{\mathbf{H}_{s}}-m^{2}-\xi R_{\mathbf{H}_{s}}\right)_{x} f(s, x, y) & =g(s, x, y)  \tag{69}\\
\text { and } \quad\left(\square_{\mathbf{H}_{s}}-m^{2}-\xi R_{\mathbf{H}_{s}}\right)_{y} f(s, x, y) & =g(s, y, x), \tag{70}
\end{align*}
$$

with some jointly smooth $g(s, x, y)$ defined on the same neighborhood of the diagonal $x=y$ as $H_{s}(x, y)$. The argument would conclude by determining a precise form of $f(s, x, y)$ near a Cauchy surface in the past of the compact region $O \subset M$, where the variation $\mathbf{H}_{s}$ differs from the reference geometry $\mathbf{H}_{0}$, and showing the existence of a unique solution of the above equations, which is moreover jointly smooth in $s, x$ and $y$ and necessarily coincides with $f(s, x, y)$.

We should note that, though the main ideas are clear, the above argument features some technical difficulties that it would take a separate paper to fully explore. For instance, it is not immediately clear which result from PDE theory would assure the existence, uniqueness and smooth parameter dependence of the solutions of Eqs. (69) and (70), all rather delicate questions, especially
in the context of partial rather than ordinary differential equations. In fact even establishing the joint smoothness of $g(s, x, y)$ goes beyond elementary facts about Hadamard parametrices (though some relevant arguments were already provided in [14, Prop. 4.1] and [28, App. A]). Also, even if the preceding issues are resolved by a clever use of standards results for hyperbolic PDEs, the domain on which Eqs. (69) and (70) are defined may not be globally hyperbolic in any meaningful way, because $g(s, x, y)$ would be defined only on a neighborhood of the diagonal in $M \times M$.

Thus, we leave the investigation of the above generalized parametrized microlocal spectral condition, of more general types of fields, of Wick powers with derivatives and of time-ordered products for future work.

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## A Peetre-Slovák's theorem with parameters

Below, we first make some remarks about how the weak regularity hypothesis (Definition 2.2) in Proposition 2.2 can be justified, despite the stronger regularity hypothesis that is usually required $[23,16,17]$. Then, we state a more general version of Proposition 2.2, in which the notion of locality is generalized to accommodate parameters. The usually complicated way in which this more general locality condition is stated is clarified through examples.

The paper [17] gives an excellent, self-contained and straight-forward proof of the version of the Peetre-Slovák theorem [17, Thm. 3.1] that we state in Proposition 2.2, with the exception that it requires the stronger regularity instead of the weak regularity hypothesis. There, the regularity hypothesis is used in exactly two places: (a) It is used once directly in the proof of Theorem 3.1, to show that the map $D[\phi]=d \circ j^{k} \phi$ factors through a smooth map $d$ on the space of $k$-jets. In that instance, the smooth "universal family of $k$-jets", which establishes the smoothness of $d$ all at once, can easily be replaced by a related smooth compactly supported variation that establishes the smoothness of $d$ in a compact neighborhood of any point of its domain. However, since that can be done for any point in the domain of $d$, the desired global smoothness of $d$ is immediate. (b) Regularity is used once more in the proof of the intermediate Lemma 2.4, to show that $D[\phi]=d \circ j^{k} \phi$ factors through some finite order jet space. There, regularity is called upon when $D$ acts on smooth families of sections $f$ and $h$ that have been constructed to have controlled behavior on a compact set $K$ by appealing to the Whitney extension theorem [17, Sec. 1.2], where $K=\left\{z_{k}\right\}_{k=0}^{\infty}$ consists of the points of a convergent sequence. Note that proof of that Lemma 2.4 goes through even with only weak regularity, provided that $f$ and $h$ could be constructed as smooth compactly supported variations, rather than smooth families. The Whitney extension theorem, as stated in [17, Sec. 1.2], which constructs smooth extensions of consistently specified jet data on
an arbitrary compact set, can only produce families of sections rather than compactly supported variations. However, it is well known that Whitney's extension theorem can be strengthened [4] to construct smooth extensions of consistently specified jet data on arbitrary closed rather than just compact sets. Using this strengthened version, it is easily seen that the above mentioned $f$ and $h$ can be constructed as smooth compactly supported variations with specified behavior both on the compact set $K$ and on the complement of any open neighborhood of $K$ with compact closure, thus showing that the proof of [17, Lem. 2.4] can be completed with only weak regularity.

The Peetre-Slovák theorem stated in Proposition 2.2 may be made significantly stronger by generalizing the admissible notion of locality. Let us now introduce the language needed to state the stronger version in a precise form. In the following, $\sigma: E \rightarrow N$ and $\rho: F \rightarrow M$, are two smooth bundles, where we have explicitly written the canonical projections, and we consider a map $D: \Gamma(E \rightarrow N) \rightarrow \Gamma(F \rightarrow M)$ between smooth sections of these bundles. We intend here to give a precise mathematical meaning to the statement that $D$ is local. Before defining the most general version of locality (cf. [16, § 18.16]), we consider several motivating cases of increasing complexity.

Case $N=M$. We say that $D$ is local when the value $\phi(x)$, for $\phi=D[\psi] \in$ $\Gamma(F \rightarrow M)$, depends only on the germ of $\psi \in \Gamma(E \rightarrow M)$ at $x \in M$. This version of locality is already sufficient for Propositions 2.1 and 2.2. We can loosen this notion of locality in several ways.

Case $N \neq M$. We may agree that $\phi(x)$, for $\phi=D[\psi] \in \Gamma(F \rightarrow M)$ and $x \in M$, may depend only on the germ of $\psi \in \Gamma(E \rightarrow N)$ at $y \in N$, with some fixed relationship $y=\chi(x)$, where $\chi: M \rightarrow N$ is some diffeomorphism. We then say that $D$ is $\chi$-local.

Case $N \neq M$ and $D$ depends on external parameters. We can introduce a bundle $\pi: P \rightarrow M$, where the manifold $P$ is interpreted as " $M$ with parameters." Then, allowing $D$ to depend on parameters means that $D$ really maps sections of $E \rightarrow N$ to sections of the pullback bundle $\pi^{*} F \rightarrow P$, interpreted as " $F$ with parameters." Let us briefly recall that, given a bundle $F \rightarrow M$ and a map $\pi: P \rightarrow M$, the pullback bundle $\pi^{*} F \rightarrow P$ is uniquely defined by the existence of a bundle morphism $\tilde{\pi}: \pi^{*} F \rightarrow F$ that is a fiber-wise isomorphism and that makes the following diagram commute


Pre-composing a section of $\pi^{*} F \rightarrow P$ with a section of $P \rightarrow M$ then yields a section of $F \rightarrow M$ given by a particular choice of parameters. Denoting $\eta=\chi \circ \pi$, we call the map $D: \Gamma(E \rightarrow N) \rightarrow \Gamma\left(\pi^{*} F \rightarrow P\right) \eta$-local when $\phi(x, p)=D[\psi](x, p)$, with $(x, p) \in P$ and $\pi(x, p)=x \in M$, depends only on the germ of $\psi$ at $y=\eta(x, p)=\chi(x) \in N$. Note that the total space of the bundle " $F$ with parameters" can be expressed as the fibered product $\pi^{*} F \cong F_{\rho} \times{ }_{\pi} P$ over $M$ (where we have explicitly named the $\rho: F \rightarrow M$ bundle projection),
which completes the pullback diagram


We can illustrate all of the above maps in the diagram

where all the solid arrows commute, the dotted arrows denote bundle sections, with $\tau: M \rightarrow P$ denoting a particular "choice of parameters," and $\phi \circ \tau$ was silently composed with the projection $F_{\rho} \times_{\pi} P \rightarrow F$.

General case. Finally, it is possible to relax the requirement that the map $\eta: P \rightarrow N$ factors as illustrated in diagram (73). The dimension of $P$ could exceed that of $N$ and $\eta$ need not be a surjection, not even a submersion. Omitting the structure of the right square of diagram (73), we also replace $F_{\rho} \times{ }_{\pi} P$ by a simple bundle $F \rightarrow P$, without requiring it to have the structure of a fibered product. So, given bundles $E \rightarrow N$ and $F \rightarrow P$, together with a smooth map $\eta: P \rightarrow N$, a map $D: \Gamma(E \rightarrow N) \rightarrow \Gamma(F \rightarrow P)$ is called $\eta$-local if $\phi(x)=D[\psi](x), x \in P$, depends only on the germ of $\psi$ at $y=\eta(x) \in N$. We can illustrate this situation by the diagram
which should be thought of as exactly analogous to diagram (73), but with the right square missing. This the rather weak notion of $\eta$-locality, with a small additional hypothesis ( $\eta$ non-locally constant), together with the condition of weak regularity (Definition 2.2) is actually sufficient for the more general version of Peetre-Slovák's theorem.

Proposition A. 1 (Peetre-Slovák's Theorem [16, § 19.10]). Let $F \rightarrow P, E \rightarrow N$ be smooth bundles and $\eta: P \rightarrow N$ a non-locally constant ${ }^{11}$ smooth map, with the interpretation as in diagram (74), and $D: \Gamma(E \rightarrow N) \rightarrow \Gamma(F \rightarrow P)$ be an $\eta$-local and weakly regular map. Then, for every compact $K \subseteq P$ and $\psi \in \Gamma(E \rightarrow N)$, there exists an integer $r$, an open neighborhood $U \subseteq J^{r}(E \rightarrow N)$ of $j^{r} \psi(N) \subset U$, with $U_{K} \subseteq U$ the subset projecting onto $\eta(K)$, and a function $d: U_{K} \rightarrow F$ that

[^9]commutes with all the projections, as illustrated by the diagram

such that $D[\xi](x)=d \circ j^{r} \xi(x)$ for any $\xi \in \Gamma(E \rightarrow N)$ with $j^{r} \xi(N) \subset U$. In other words, $D$ is a differential operator of locally finite order, where locality is with respect to compact subsets of $P$ and compact open neighborhoods in $\Gamma(E \rightarrow N)$.
Sketch of proof. With the definitions as discussed above, the proposition is essentially a restatement of Theorem 19.10 of [16], which follows directly from Theorem 19.7 and Corollary 19.8 that precede it. We refer the reader to the book [16] for full details. Let us simply mention that, in general outline, the proof proceeds by contradiction. If $D$ depended non-trivially on an infinite number of derivatives of its argument, then it would be possible to engineer a smooth section $\psi$ such that $D[\psi]$ could not itself be smooth. While the proofs in [16] rely on regularity instead of weak regularity, the weaker hypothesis is actually sufficient, as we have discussed earlier.

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[^1]:    ${ }^{1}$ As in [13], we will always treat $m^{2}$ as a real number, which could be either positive, zero, or even negative, as ultraviolet renormalization is not sensitive to the sign of $m^{2}$.

[^2]:    ${ }^{2}$ On page 490 of [2], the unpublished report [1] is used as the main reference for the properties of these coordinates. In [26] it is explained further that their origin goes back to at least [21] and even earlier to [24].

[^3]:    ${ }^{3}$ In an earlier version of the manuscript, we mistakenly omitted the regularity hypothesis from the statement of the Peetre-Slovák theorem. We thank the anonymous referee for bringing that to our attention.

[^4]:    ${ }^{4}$ The homogeneous degree of $P(g, \varepsilon(g))$ counts the components of $g$ with degree 2 and the components of $\varepsilon(g)$ with degree $n$.

[^5]:    ${ }^{5}$ For every Hadamard quasifree state $\omega$ over $W(M, \mathbf{g})$ the map $\mathcal{D}(M) \ni f \mapsto \omega(\Phi(f))$ is a distribution in the proper sense. A weaker requirement allowing to smear fields with distributions of a suitable wavefront set can be given exploiting the so called Hörmander pseudotopology [13], but it is irrelevant for this work.
    ${ }^{6}$ Note that $\mathcal{D}: \mathfrak{M a n} \rightarrow \mathfrak{L C V}$ is a (covariant) functor from manifolds to locally convex topological vector spaces. It assigns the space of complex valued test functions $\mathcal{D}(M)$ to a manifold $M$ and maps a morphism $\chi: M \rightarrow M^{\prime}$ to the induced extension by zero map $\chi_{*}: \mathcal{D}(M) \rightarrow \mathcal{D}\left(M^{\prime}\right)$.

[^6]:    ${ }^{7}$ We make the tacit and harmless assumption that we are working with small categories, whose objects and morphisms constitute sets.

[^7]:    ${ }^{8}$ The differences between the Hollands and Wald definition of Wick powers from ours is detailed in Remark 3.3.

[^8]:    ${ }^{9}$ In principle, the hypothesis of locally bounded order tells us that this is true for a sufficiently small neighborhood of $y$, with $r$ possibly increasing on larger neighborhoods. However, since such a neighborhood exists around any $y \in U$, a simple argument based on open covers and the compactness of $U$ shows that the order $r$ can be chosen uniformly over an arbitrary compact $U$.

[^9]:    ${ }^{11}$ By non-locally constant we mean that for every open $U \subseteq P$ the image $\eta(U)$ contains at least two points.

