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# A Note on the Inverse Scattering Problem in Quantum Field Theory 

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#### Abstract

We study the set of local fields $\phi$ describing the dynamics of a scalar, massless particle. It turns out that these fields are relatively local to the free, massless, scalar field $A$ if the massless particle does not interact. This leads to a simple algebraic characterisation of interacting fields in the above framework.


## 1. Problems and Results

An old problem in quantum field theory is to characterize all local fields leading to a given $S$-matrix [1], [2; Chapter 4.6], [3; Chapter 19.5]. Our interest in this question arose from a desire to have a local criterion distinguishing field theories with interaction from non-interacting ones. For that purpose it would be sufficient to know all fields leading to a trivial $S$-matrix. However, even this simpler problem has not yet been solved.

In view of this situation we found it worth while putting the following analysis on record although it applies only to a rather special case. We consider in this paper the model of a scalar, massless particle which does not interact. It turns out that there each interpolating field $\phi$ for the (trivial) $S$-matrix is an element of the Borchers class [1] of the free, massless field $A$. All these elements are explicitly known (see [4] and the Appendix) and this solves the inverse scattering problem for this special model. Moreover, because the free, massless field commutes with itself at timelike separations, one obtains a simple algebraic characterisation of interacting fields in the above framework : a local field $\phi$, describing the dynamics of a scalar, massless particle leads to a non-trivial $S$-matrix if and only if the commutator $[\phi(x), \phi(y)]$ does not vanish at timelike distances $(x-y)$.

Our argument is based on the following reasoning: if the asymptotic fields coincide, $\phi^{\text {in }}=\phi^{\text {out }}=A$, they have the same TCP-operator as the interpolating field $\phi$; consequently, $\phi$ and $A$ are weakly local with respect to each other [1]. Moreover, by Huyghens' principle [5] the commutator between $\phi(x)$ and $A(y)$ vanishes at timelike distances $(x-y)$. Hence if $B$ is any element of the Borchers class of $A$ the support of the vacuum expectation value

$$
\begin{equation*}
K(x)=(\Omega,[B(0), \phi(x)] \Omega) \tag{1}
\end{equation*}
$$

is confined to the lightcone $x^{2}=0$. It follows then from the spectrum condition and temperateness that

$$
\begin{equation*}
K(x)=\sum_{\text {finite }} c_{(m), n} x^{(m)} \varepsilon\left(x_{0}\right) \delta^{(n)}\left(x^{2}\right), \tag{2}
\end{equation*}
$$

where $(m)$ is a multiindex and $x^{(m)}=x_{0}^{m_{0}} x_{1}^{m_{1}} x_{2}^{m_{2}} x_{3}^{m_{3}}$. Since the right hand side of this equation is a sum of homogeneous distributions we conclude that the unitary operators $D(\lambda), \lambda>0$ which induce the dilations of the asymptotic field $A$ act on $\phi$ according to

$$
\begin{equation*}
D(\lambda) \phi\left(\lambda^{-1} x\right) D(\lambda)^{-1}=\sum_{d=0}^{n} \lambda^{d} \phi_{d}(x), \tag{3}
\end{equation*}
$$

i.e. $\phi$ is a finite sum of fields carrying a dimension ; in particular $\phi_{0}=(\Omega, \phi(0) \Omega) \cdot 1$ and $\phi_{1}=A$. Locality of $\phi$ then leads to recursive relations for $\phi_{d}$ which imply that these fields are relatively local to $A$.

Admittedly, this method of proof is tailored to the massless case and cannot be generalized to massive models. Nevertheless we hope that this contribution will serve to stimulate further investigations of these problems.

## 2. Details

The assumption that we are dealing with a field $\phi$ which leads to a trivial $S$-matrix can be expressed as follows : let $A$ be the free, massless, scalar field acting in Fockspace $\mathscr{H}$. We identify $A$ with the incoming field constructed from the real, local, temperate field $\phi$, i.e. we require
a) $\phi$ transforms as a scalar field under the same unitary representation $(\Lambda, x) \rightarrow U(x) \cdot U(\Lambda)$ of the Poincaré group as $A$.
b) $P_{1} \phi(x) \Omega=A(x) \Omega$ where $P_{1}$ is the projection onto the one-particle states in $\mathscr{H}$.
c) $[A(x), \phi(y)]=0$ for $(x-y) \in V_{+}$(Huyghens' principle [5]).

These assumptions characterize $A$ unambiguously as the incoming field constructed from $\phi^{1}$. If $A$ were the outgoing field the only change would be in relation c) where we would have to replace $V_{+}$by the backward cone $V_{-}$. Hence, if there is no scattering in the model, relation c) holds for $(x-y) \in V_{+} \cup V_{-}$. An alternative way of expressing the triviality of the $S$-matrix is [1]
d) $\phi$ is weakly local relative to $A$.

For completeness let us also specify the domains of $A$ and $\phi$. We suppose that we may apply arbitrary smeared polynomials in $A$ and $\phi$ to the vacuum $\Omega$. This rather strong assumption is not really necessary for the proof. However, it will allow us to neglect intricate domain questions in what follows.

Now let $D(\lambda), \lambda>0$ be the dilation operators acting on $A$ according to

$$
\begin{equation*}
D(\lambda) A\left(\lambda^{-1} x\right) D(\lambda)^{-1}=\lambda \cdot A(x) . \tag{4}
\end{equation*}
$$

[^0]Since the Borchers class $\mathfrak{B}(A)$ of $A$ contains only Wick-polynomials of $A$ and its derivatives (see the Appendix) each $B \in \mathfrak{B}(A)$ can be decomposed into a finite sum of fields $B_{d} \in \mathfrak{B}(A)$ with dimension $d \in(0, \mathbb{N})$. We take any such $B_{d}$ and analyse the distribution

$$
\begin{equation*}
K_{\lambda}(x)=\left(\Omega,\left[B_{d}(0), \phi_{\lambda}(x)\right] \Omega\right)=\lambda^{-d}\left(\Omega,\left[B_{d}(0), \phi\left(\lambda^{-1} x\right)\right] \Omega\right), \tag{5}
\end{equation*}
$$

where $\phi_{\lambda}(x)=D(\lambda) \phi\left(\lambda^{-1} x\right) D(\lambda)^{-1}$. It follows from assumptions c) and d) that $\left(\Omega,\left[B_{d}(0), \phi(x)\right] \Omega\right)$ has support on the light cone $x^{2}=0$. The Fourier transform of this distribution vanishes for $p^{2}<0$ owing to the spectrum condition. We give the most general form of such a distribution in the subsequent lemma.

Lemma 1. Let $S \in \mathscr{S}^{\prime}\left(\mathbb{R}^{4}\right)$ be a temperate distribution with $\operatorname{supp} S \subseteq\left\{x^{2}=0\right\}$ and $\operatorname{supp} \tilde{S} \cong\left\{p^{2} \geqq 0\right\}$. Then

$$
S(x)=\sum_{\text {finite }} c_{(m), n} x^{(m)} \varepsilon\left(x_{0}\right) \delta^{(n)}\left(x^{2}\right)
$$

where $x^{(m)}=x_{0}^{m_{0}} x_{1}^{m_{1}} x_{2}^{m_{2}} x_{3}^{m_{3}}$ with $m_{i} \in(0, \mathbb{N}) . \delta^{(n)}(\cdot)$ denotes the $n$-th derivative of the $\delta$ function and $\varepsilon(\cdot)$ the sign-function.

Proof. Since $S$ has its support on the surface $x^{2}=0$ and is temperate there exists a minimal number $N \in \mathbb{N}$ such that $\left(x^{2}\right)^{N} S(x)=0$, hence $\square_{p}^{N} \tilde{S}(p)=0$. We divide the set of distributions $S$ into subsets corresponding to $N$ and prove the lemma by induction in $N$ : for $N=1 \tilde{S}(p)$ is a solution of the wave equation and can be expressed by its (temperate) Cauchy data on the hyperplane $p_{0}=0$. Because of the fact that supp $\tilde{S} \subseteq\left\{p^{2} \geqq 0\right\}$, these data are localized at $\boldsymbol{p}=0$ and therefore finite sums of $\delta$-functions and their derivatives. Taking into account that $\varepsilon\left(p_{0}\right) \delta\left(p^{2}\right)$ is a solution of the wave equation with Cauchy data 0 and $2 \pi \cdot \delta(\boldsymbol{p})$ it is then easy to verify that

$$
\tilde{S}(p)=P_{1}(\partial) \varepsilon\left(p_{0}\right) \delta\left(p^{2}\right)+P_{2}(\partial) \partial_{0} \varepsilon\left(p_{0}\right) \delta\left(p^{2}\right),
$$

where $P_{1}, P_{2}$ are certain polynomials and $\left(\partial_{0}, \boldsymbol{\partial}\right)$ are the derivatives with respect to $\left(p_{0}, \boldsymbol{p}\right)$. This proves the statement for $N=1$ after Fourier transformation. To complete the induction we must show that if $x^{2} S(x)$ is of the form given in the lemma, i.e.

$$
x^{2} \cdot S(x)=\sum_{\text {finite }} c_{(m), n} x^{(m)} \varepsilon\left(x_{0}\right) \delta^{(n)}\left(x^{2}\right),
$$

then $S(x)$ itself has this form. To this end we must divide the above equation by $x^{2}$. An obvious solution to this problem with the desired properties is

$$
S_{0}(x)=-\sum_{\text {finite }} c_{(m), n} x^{(m)}(n+1)^{-1} \varepsilon\left(x_{0}\right) \delta^{(n+1)}\left(x^{2}\right) .
$$

But the difference between $S$ and $S_{0}$ is a distribution $\Delta S$ with supp $\tilde{\Delta} S \subseteq\left\{p^{2} \geqq 0\right\}$ and $x^{2} \cdot \Delta S(x)=0$. We may therefore apply to $\Delta S$ the argument given above for $N=1$ and conclude that $S=S_{0}+\Delta S$ has the form given in the lemma.

This lemma shows that the commutator function $K_{\lambda}$ given in (5) has a finite Laurent expansion at $\lambda=0$,

$$
\begin{equation*}
K_{\lambda}(x)=\sum_{\text {finite }} \lambda^{i} K_{i}(x) \quad \text { with } \quad i \in \mathbb{Z} \tag{6}
\end{equation*}
$$

In fact $K_{\lambda}$ is even a polynomial in $\lambda$; the degree of this polynomial is bounded by a number $n$ which depends on $\phi$ but not on $B_{d}$. To verify this we estimate

$$
\begin{equation*}
\left|\int d^{4} x d^{4} y g(x) h(y) K_{\lambda}(x-y)\right| \leqq\left\|\phi_{\lambda}(g) \Omega\right\| \cdot\left\{\left\|B_{d}(h) \Omega\right\|+\left\|B_{d}(h)^{*} \Omega\right\|\right\} \tag{7}
\end{equation*}
$$

with $g$, $h$ real. Now $(\Omega, \phi(0) \phi(x) \Omega)$ is the Fourier transform of a temperate measure and therefore $\left\|\phi_{\lambda}(g) \Omega\right\| \leqq c_{g}\left(1+\lambda^{n}\right)$ for some $n \in \mathbb{N}$. Thus $\left|K_{\lambda}(f)\right| \leqq c_{f}^{\prime}\left(1+\lambda^{n}\right)$ and this limits the exponents $i$ in relation (6) to $0 \leqq i \leqq n$. Taking into account the spectrum condition we conclude that for arbitrary $B_{d} \in \mathfrak{B}(A)$

$$
\begin{equation*}
\frac{d^{n+1}}{d \lambda^{n+1}}\left(\Omega, B_{d}(0) \phi_{\lambda}(x) \Omega\right)=0 \tag{8}
\end{equation*}
$$

But the linear span of $\left\{B_{d}(f) \Omega: B_{d} \in \mathfrak{B}(A), d \in(0, \mathbb{N}), f \in \mathscr{S}\left(\mathbb{R}^{4}\right)\right\}$ is dense in $\mathscr{H}^{2}$ and therefore

$$
\begin{equation*}
\frac{d^{n+1}}{d \lambda^{n+1}} \phi_{\lambda}(x) \Omega=0 \tag{9}
\end{equation*}
$$

This establishes relation (3) on the vacuum. Now it follows from Huyghens' principle that $\left[\phi_{\lambda}(x), A(y)\right]=0$ for $(x-y)^{2}>0$ and all $\lambda>0$. So relation (9) still holds if we replace $\Omega$ by a vector $\mathfrak{P}(A) \Omega$, where $\mathfrak{P}(A)$ is any smeared polynomial in $A$. This establishes relation (3) on the dense domain $\mathscr{D}_{0}=\{\mathfrak{P}(A) \Omega\}$.

In the second step of our argument we analyse the components $\phi_{d}$ in the decomposition $\phi_{\lambda}=\sum_{d=0}^{n} \lambda^{d} \phi_{d}$. Clearly the $\phi_{d}$ are scalar fields which are defined on $\mathscr{D}_{0} ;$ moreover they satisfy $\left[\phi_{d}(x), A(y)\right]=0$ for $(x-y)^{2}>0$. Keeping this in mind and using assumption b ) it is easy to see that $\phi_{0}=(\Omega, \phi(0) \Omega) \cdot 1$ and $\phi_{1}=A$ [7]. In order to show that $\phi_{d} \in \mathfrak{B}(A)$ for $d \geqq 2$ we start from the relation $\left[\phi_{\lambda}(x), \phi_{\lambda}(y)\right]=0$ which holds for $(x-y)^{2}<0$ because of locality. Then we replace $\phi_{\lambda}$ by the expression given above and arrive at the equations $(1 \leqq m \leqq n)$

$$
\begin{equation*}
\sum_{d=1}^{m}\left[\phi_{d}(x), \phi_{m+1-d}(y)\right]=0 \quad \text { for } \quad(x-y)^{2}<0 \tag{10}
\end{equation*}
$$

which hold in the sense of bilinear forms on $\overline{\mathscr{D}_{0}} \times \mathscr{D}_{0}$. But $\phi_{1}=A$ with $\square A(y)=0$ and therefore $(3 \leqq m \leqq n)$

$$
\begin{equation*}
\left[A(x), \square \phi_{m}(y)\right]=-\sum_{d=2}^{m-1}\left[\phi_{d}(x), \square \phi_{m+1-d}(y)\right] \quad \text { for } \quad(x-y)^{2}<0 \tag{11}
\end{equation*}
$$

[^1]Hence if $\phi_{2}, \ldots \phi_{m-1} \in \mathfrak{B}(A)$, the right hand side of this equation vanishes and $\square \phi_{m}$ must be an element of $\mathfrak{B}(A)$. To complete the argument it suffices therefore to show that $\square \phi_{d} \in \mathfrak{B}(A)$ implies $\phi_{d} \in \mathfrak{B}(A)$ for $d \geqq 2$.

Lemma 2. Let $\phi_{d}, d \geqq 2$ be a temperate field with properties specified above. If $\square \phi_{d} \in \mathfrak{B}(A)$ then also $\phi_{d} \in \mathfrak{B}(A)$.
Proof. Since $\square \phi_{d} \in \mathfrak{B}(A)$ it can be represented in the form

$$
\begin{equation*}
\square \phi_{d}(x)=\sum_{k=2}^{d+2} P_{k}\left(\partial_{1}, \ldots \partial_{k}\right): A\left(x_{1}\right) \ldots A\left(x_{k}\right):\left.\right|_{x_{1}=\ldots=x_{k}=x} \tag{*}
\end{equation*}
$$

where $P_{k}$ are symmetric polynomials in the derivatives $\partial_{i}=\left(\partial_{i_{0}}, \partial_{i}\right)$ with respect to $\left(x_{i_{0}}, \boldsymbol{x}_{i}\right)$. We shall see that each $P_{k}\left(\partial_{1}, \ldots \partial_{k}\right)$ contains a factor $\square_{k}=\left(\partial_{1}+\ldots+\partial_{k}\right)^{2}$. It is then obvious that $\phi_{d}$ itself has the form $(*)$ up to a term which is a solution of the "homogeneous equation" $\square \phi_{d}(x)=0$. However, this equation has only trivial solutions if $d \geqq 2$ since then $\phi_{d}(x) \Omega$ has no one-particle contribution and the vacuum is separating for the operators $\phi_{d}$.

The assertion concerning the polynomials $P_{k}$ will be proved by induction. For $k=2$ the statement follows simply from Lorentz-invariance and the fact that $\square \phi_{d}$ has dimension $d+2$; so $P_{2}\left(\partial_{1}, \partial_{2}\right)$ must be a homogeneous polynomial of degree $d \geqq 2$ in the invariant $\left(\partial_{1}+\partial_{2}\right)^{2}$. (The invariants $\partial_{1}^{2}, \partial_{2}^{2}$ do not appear because of $\square A(x)=0$.)

Now if each $P_{k}\left(\partial_{1}, \ldots \partial_{k}\right), 2 \leqq k \leqq l-1$ contains a factor $\square_{k}$ we may subtract in relation $(*)$ the corresponding contributions from $\square \phi_{d}$ and assume that the sum starts with $l \geqq 3$. Then the expressions $\int d^{4} x \square \phi_{d}(x) D^{\binom{\text {ret }}{\mathrm{adv}}}(x)$ and $\int d^{4} x \square \phi_{d}(x) D(x)$ are well defined as bilinear forms on $\overline{\mathscr{D}}_{0} \times \mathscr{D}_{0}$. (Here $D$ denotes the zero mass PauliJordan commutator function and $D^{\binom{\text {ret }}{\text { adv }}}$ the retarded and advanced parts of $D$.) Moreover we have $\int d^{4} x \square \phi_{d}(x) D(x)=0$ on $\overline{\mathscr{D}}_{0} \times \mathscr{D}_{0}$ as expected from the YangFeldman equations. To confirm this statement we consider the temperate distribution

$$
\begin{equation*}
\int d^{4} x D(x)\left(\Omega, A\left(x_{1}\right) \ldots A\left(x_{k}\right) \square \phi_{d}(x) A\left(x_{k+1}\right) \ldots A\left(x_{l}\right) \Omega\right) \tag{**}
\end{equation*}
$$

Taking into account the support properties of $D^{\binom{\text {(ret }}{\mathrm{adv}}}, D=D^{(\mathrm{ret})}-D^{\text {(adv) }}$ and the timelike commutation relations between $\phi_{d}$ and $A$ it follows ${ }^{3}$ that ( $* *$ ) coincides with

$$
\begin{aligned}
& \int d^{4} x D^{(\mathrm{ret})}(x)\left(\Omega, \square \phi_{d}(x) A\left(x_{1}\right) \ldots A\left(x_{l}\right) \Omega\right) \\
& -\int d^{4} x D^{(\text {adv })}(x)\left(\Omega, A\left(x_{1}\right) \ldots A\left(x_{l}\right) \square \phi_{d}(x) \Omega\right),
\end{aligned}
$$

provided $x_{1}, \ldots x_{k} \in V_{-}$and $x_{k+1}, \ldots x_{l} \in V_{+}$. But

$$
\int d^{4} x D^{(\mathrm{ret})}(x) \square \phi_{d}(x) \Omega=\int d^{4} x D^{(\text {adv })}(x) \square \phi_{d}(x) \Omega=\phi_{d}(0) \Omega
$$

if $d \geqq 2$ and using the timelike commutation relations once more it becomes clear that $(* *)$ vanishes for the special configurations $x_{1}, \ldots x_{l}$ mentioned above. On the

[^2]other hand (**) is the boundary value of a function which is analytic in a tube because of the spectrum condition. We may therefore apply the "edge of the wedge" theorem [2] and conclude that ( $* *$ ) vanishes identically. If we insert now into ( $* *$ ) the expression given for $\square \phi_{d}$ in relation $(*)$ we get after a simple calculation
$$
\int \prod_{j=1}^{l} d^{4} p_{j} \delta_{+}\left(p_{j}^{2}\right) \delta\left(\left(\sum_{j=1}^{l} s_{j} p_{j}\right)^{2}\right) \varepsilon\left(\sum_{j=1}^{l} s_{j} p_{j_{0}}\right) \cdot P_{l}\left(i s_{1} p_{1}, \ldots i s_{l} p_{l}\right) \exp \left(i \sum_{j=1}^{l} s_{j}\left(p_{j} x_{j}\right)\right)=0
$$
where $s_{j}=-1$ for $j \leqq k$ and $s_{j}=1$ for $j \geqq k+1$. Since the measure
$$
\prod_{j=1}^{l} d^{4} p_{j} \delta_{+}\left(p_{j}^{2}\right) \delta\left(\left(\sum_{j=1}^{l} s_{j} p_{j}\right)^{2}\right)
$$
has support on the manifold $\left(\sum_{j=1}^{l} s_{j} p_{j}\right)^{2}=0, p_{j}^{2}=0, p_{j_{0}} \geqq 0, j=1, \ldots l$ if $1 \leqq k \leqq l-1$ and $l \geqq 3$, it follows from the above equation that $P_{l}\left(i s_{1} p_{1}, \ldots i s_{l} p_{l}\right)$ vanishes at these points. Thus $P_{l}\left(i p_{1}, \ldots i p_{l}\right)$ vanishes in particular in an open neighbourhood of a regular point on the manifold
$$
p_{1}^{2}=\ldots=p_{l}^{2}=0, \quad\left(\sum_{j=1}^{l} p_{j}\right)^{2}=0
$$

Each polynomial with this property is of the form

$$
\left(\sum_{j=1}^{l} p_{j}\right)^{2} \cdot Q\left(p_{1}, \ldots p_{l}\right)+\sum_{j=1}^{l}\left(p_{j}\right)^{2} \cdot R_{j}\left(p_{1}, \ldots p_{l}\right)
$$

where $Q, R_{j}$ are also polynomials (see the Appendix). Again the terms containing $p_{1}^{2}, \ldots p_{l}^{2}$ do not contribute and therefore $P_{l}\left(i p_{1}, \ldots i p_{l}\right)$ contains a factor $\left(\sum_{j \geqq 1}^{l} p_{j}\right)^{2}$. This finishes the proof of the lemma.

Collecting the results in this chapter we arrive at the
Theorem. Let $\phi$ be a real, local, temperate field with properties a) to d) given above. Then $\phi \in \mathfrak{B}(A)$.

An immediate consequence of this theorem is the
Corollary. Let $\phi$ be a real, local, temperate field with properties a) to c) given above. Then $\phi$ leads to a trivial S-matrix if and only if

$$
[\phi(x), \phi(y)]=0 \quad \text { for } \quad(x-y)^{2}>0
$$

## Appendix

a) The Borchers Class of $A$

For the convenience of the reader we recall here some facts about the Borchers class $\mathfrak{B}(A)$ of the free, massless, scalar field $A$ : a temperate field $B$ is said to be an element of $\mathfrak{B}(A)$ if it has the following properties:
a) $B(f), f \in \mathscr{S}\left(\mathbb{R}^{4}\right)$ is defined on the dense set of vectors $\mathscr{D}_{0}=\{\mathfrak{P}(A) \Omega\}$ where $\mathfrak{P}(A)$ are the smeared polynomials in $A$.
b) $B$ transforms under the same unitary representation $x \rightarrow U(x)$ of the translations as $A$.
c) $[A(x), B(y)]=0$ for $(x-y)^{2}<0$ in the sense of bilinear forms on $\overline{\mathscr{D}}_{0} \times \mathscr{D}_{0}$.

Following an argument of Epstein [4] we sketch the proof that $\mathfrak{B}(A)$ consists of all Wick polynomials of $A$ and its derivatives. For this purpose we consider the $n$ fold commutator function

$$
\left(\Omega,\left[A\left(x_{1}\right),\left[A\left(x_{2}\right), \ldots\left[A\left(x_{n}\right), B(0)\right] \ldots\right] \Omega\right)\right.
$$

which is a solution of the wave equation in each of its arguments. This expression is symmetric in $x_{1}, \ldots x_{n}$ (because $A$ has $c$-number commutation relations) and therefore vanishes if any one of its arguments $x_{i}$ becomes spacelike. Hence if we express the commutator function by its Cauchy data at time $x_{1_{0}}=\ldots=x_{n_{0}}=0$ we realize that it is of the form

$$
P_{n}\left(\partial_{1}, \ldots \partial_{n}\right) \prod_{i=1}^{n} D\left(x_{i}\right)
$$

where $P_{n}$ is some symmetric polynomial in the derivatives $\partial_{i}=\left(\partial_{i_{0}}, \partial_{i}\right)$ and $D$ is the zero mass Pauli-Jordan commutator function. We get therefore in the sense of bilinear forms on $\overline{\mathscr{D}_{0}} \times \mathscr{D}_{0}$ :

$$
B(x)=\sum_{n=0}^{\infty} \frac{1}{n!} P_{n}\left(-\partial_{1}, \ldots-\partial_{n}\right): A\left(x_{1}\right) \ldots A\left(x_{n}\right):\left.\right|_{x_{1}=\ldots=x_{n}=x}
$$

But this series must terminate at some finite $n$ because $B$ is a temperate field.

## b) Structure of the Polynomials $P_{l}$

In this Appendix we want to show that a polynomial $P\left(p_{1}, \ldots p_{l}\right), l \geqq 3$ which vanishes in a real neighbourhood of a regular point on the manifold $p_{1}^{2}=\ldots=p_{l}^{2}$

$$
\begin{aligned}
= & \left(\sum_{j=1}^{l} p_{j}\right)^{2}=0 \text { is of the form } \\
& P\left(p_{1}, \ldots p_{l}\right)=\left(\sum_{j=1}^{l} p_{j}\right)^{2} \cdot Q\left(p_{1}, \ldots p_{l}\right)+\sum_{j=1}^{l} p_{j}^{2} \cdot R_{j}\left(p_{1}, \ldots p_{l}\right)
\end{aligned}
$$

with polynomials $Q$ and $R_{j}$. Unfortunately, we have not been able to give an elementary proof of this statement. We shall heavily rely on results of the theory of polynomial rings, as expounded e.g. in [8] and [9].

To begin with we mention that a polynomial $P$ with properties given above, automatically vanishes on the whole complex manifold $p_{1}^{2}=\ldots=p_{l}^{2}=\left(\sum_{j=1}^{l} p_{j}\right)^{2}=0$. This may be seen either by function theoretical techniques or from the fact that this manifold is algebraically irreducible. (It is of course crucial here that $P$ vanishes in a whole neighbourhood of a regular point on this manifold.)

In a second step we show that $P$ has the desired form up to a factor. To this end we introduce an auxiliary variable $\lambda$ and study the polynomial $P^{\prime}\left(\lambda, p_{1}, \ldots p_{l}\right)$ $=P\left(\lambda \cdot p_{1}, p_{2}, \ldots p_{l}\right) . P^{\prime}$ vanishes if $p_{1}^{2}=\ldots=p_{l}^{2}=0$ and

$$
\lambda=\lambda_{0}=-\left(\sum_{j \geqq 2}\left(\mathrm{p}_{1} \mathrm{p}_{j}\right)\right)^{-1} \cdot\left(\sum_{j>k \geqq 2}\left(\mathrm{p}_{j} \mathrm{p}_{k}\right)\right) .
$$

Hence if $m$ is the degree of $P^{\prime}$ in the variable $\lambda$ it follows that

$$
R^{\prime}\left(p_{1}, \ldots p_{l}\right)=\left(\sum_{j=2}^{l}\left(p_{1} p_{j}\right)\right)^{m} \cdot P^{\prime}\left(\lambda_{0}, p_{1}, \ldots p_{l}\right)
$$

is a polynomial which vanishes if $p_{1}^{2}=\ldots=p_{l}^{2}=0$. A straightforward application of the Euclidean division algorithm then gives

$$
R^{\prime}\left(p_{1}, \ldots p_{l}\right)=\sum_{j=1}^{l} p_{j}^{2} \cdot R_{j}^{\prime}\left(p_{1}, \ldots p_{l}\right)
$$

with polynomials $R_{j}^{\prime}$. Since $P^{\prime}(\lambda=1)=P$ it is also obvious that $\left[\left(\sum_{j=2}^{l}\left(p_{1} p_{j}\right)\right)^{m} \cdot P-R^{\prime}\right]$ can be divided by $\sum_{j>k \geqq 1}\left(p_{j} p_{k}\right)$. Therefore we get

$$
\begin{equation*}
N \cdot P=\left(\sum_{j=1}^{l} p_{j}\right)^{2} \cdot Q+\sum_{j=1}^{l} p_{j}^{2} \cdot R_{j} \tag{*}
\end{equation*}
$$

with polynomials $Q, R_{j}$, and $N=\left(\sum_{j=2}^{l}\left(p_{1} p_{j}\right)\right)^{m}$.
The last (and non-trivial) step in our analysis consists of showing that we may choose the polynomials $Q$ and $R_{j}$ in the above expression in such a way that they can be divided by $N$. For this purpose we introduce some algebraic notions.

Let $f_{1}, \ldots f_{k}$ be fixed polynomials. We call the set of all polynomials of the form ' $g_{1} f_{1}+\ldots+g_{k} f_{k}$ ), $g_{i}$ being arbitrary polynomials, the ideal $\mathfrak{I}=\left(f_{1}, \ldots f_{k}\right)$ generated by $f_{1}, \ldots f_{k} \cdot \mathfrak{I}$ is called primary if $f \cdot g \in \mathfrak{I}$ and $f \notin \mathfrak{I}$ implies $g^{m} \in \mathfrak{I}$ for some $m \in \mathbb{N}$. It is one of the fundamental results in the theory of polynomial rings that each ideal $\mathfrak{I}$ is a finite intersection of primary ideals, $\mathfrak{J}=\bigcap_{i=1}^{m} \mathfrak{I}_{i}$. Such a representation is of course not unique. Let us pick a representation in which one cannot omit any $\mathfrak{I}_{i}$ and consider the varieties of the primary ideals $\mathfrak{I}_{i}$, i.e. the set of all common zeros of the polynomials in $\mathfrak{J}_{i}$. It turns out that these varieties are uniquely determined by $\mathfrak{I}$. So let us call them associated varieties. Now it is crucial for our argument that if $f$ is a polynomial, which does not identically vanish on any of these varieties and if $f \cdot g \in \mathfrak{I}$ then $g$ must be an element of $\mathfrak{J}$ [8, Chapter 16].

So in our case we have only to show that $N$ does not identically vanish on any of the associated varieties of the ideal $\mathfrak{J}_{0}=\left(p_{1}^{2}, \ldots p_{1}^{2},\left(\sum_{j=1}^{l} p_{j}\right)^{2}\right)$. To this end we must determine them. First of all we realize that the associated varieties of $\mathfrak{I}_{0}$ are subsets of the variety of $\mathfrak{I}_{0}$, i.e. the manifold $\mathfrak{B}_{0}=\left\{p_{1}^{2}=\ldots=p_{l}^{2}=\left(\sum_{j=1}^{l} p_{j}\right)^{2}=0\right\}$; therefore
their dimensions must be less or equal to $(3 l-1)$. But then Macaulay's theorem [ 9 , Theorem 26] tells us that each of these varieties has dimension $(3 l-1)$, because $\mathfrak{I}_{0}$ is generated by $4 l-(3 l-1)$ functions. Bearing in mind that the manifold $\mathfrak{B}_{0}$ is algebraically irreducible we conclude that the associated varieties of $\mathfrak{J}_{0}$ coincide with $\mathfrak{B}_{0}$. Now we come back to relation $(*)$ : since $N \cdot P \in \mathfrak{I}_{0}$ and $N$ does not identically vanish on $\mathfrak{B}_{0}$ if $l \geqq 3, P$ must be an element of $\mathfrak{I}_{0}$. This, finally, proves the statement.

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[^0]:    It is of course a more delicate problem whether, given a $\phi$, one can always construct an incoming ield $A$ with properties specified above. This question has been answered affirmatively in the algebraic ramework of field theory [5]

[^1]:    2 This may be verified by direct computation or extracted from [6]

[^2]:    3 Our argument is not completely rigorous in view of the distribution character of $D^{\binom{\mathrm{ret})}{\mathrm{adv}}}$ and $D$. However, it can be straightened out if one approximates these distributions by suitable testfunctions $D_{n}^{\binom{\text {ret }}{\text { adv }}}, D_{n}$ and takes limits afterwards

