# Variational Principle for Derivation of Macroscopic Equations 

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Received June 18, 1980


#### Abstract

It is pointed out that the Schwinger variational principle of scattering theory applies to the case of linear and nonlinear relaxation problems in quantum statistics. By means of this principle it is possible to derive closed sets of equations for expectation values. To illustrate this variational method and to clarify the connection to other standard approaches some simple examples are treated for which the equations of motion are already known.


KEY WORDS: Variational principle; von Neumann equation; equations for macroscopic observables; relaxation.

## 1. INTRODUCTION

For classical dynamic correlation functions of linear response a Schwingertype variational principle has been formulated. ${ }^{(1)}$ Later the same idea of such a variational principle has independently been introduced for quantum systems. ${ }^{(2)}$ It is the purpose of this paper to point out that this variational principle, formally taken from scattering theory, can be extended to the case of general dynamics of macroscopic observables. In this general form the variational principle can be used to derive closed systems of equations for macroscopic variables from the microscopic von Neumann equation. In this paper we will give no new equations of motion. We wish only to establish the variational principle and to demonstrate how it works. To this end we treat simple physical examples and show how standard equations of statistical mechanics emerge from this variational approach. The examples considered can be a guide to find a suitable variational ansatz in complicated cases.

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## 2. STATIONARY VARIATIONAL PRINCIPLE

Consider a system with time-independent Hamiltonian $\mathscr{H}$ described by a statistical operator $\rho$. Then the time evolution of any observable $A_{v}$ can be expressed as follows:

$$
\begin{equation*}
\left\langle A_{v}\right\rangle(t)=\operatorname{Tr} \rho e^{i L t} A_{v} \tag{1}
\end{equation*}
$$

where the Liouvillian $L$ is defined by

$$
\begin{equation*}
L A=(1 / \hbar)[\mathscr{K}, A] \tag{2}
\end{equation*}
$$

Now from (1) it is evident that the Laplace transform of $\left\langle A_{v}\right\rangle(t)$ can be written as a resolvent

$$
\begin{equation*}
\left\langle A_{v}\right\rangle(z)=-i \int_{0}^{\infty} d t e^{i z t}\left\langle A_{v}\right\rangle(t)=\operatorname{Tr} \rho(z+L)^{-1} A_{v} \tag{3}
\end{equation*}
$$

The calculation of a resolvent is a basic mathematical problem in scattering theory. ${ }^{(3)}$ Hence we can transfer the methods used there to quantum statistics and write down a Schwinger-type functional

$$
\begin{equation*}
F\left[\phi^{\prime}, \phi_{v}\right]=\operatorname{Tr} \phi^{\prime} A_{v}+\operatorname{Tr} \rho \phi_{v}-\operatorname{Tr} \phi^{\prime}(z+L) \phi_{v} \tag{4}
\end{equation*}
$$

which has to be varied with respect to $\phi^{\prime}$ and $\phi_{v}$. It is a simple matter to show that $F$ is stationary for

$$
\begin{align*}
\phi_{v \mathrm{st}} & =(z+L)^{-1} A_{v}=A_{v}(z)  \tag{5a}\\
\phi_{\mathrm{st}}^{\prime} & =(z-L)^{-1} \rho=\rho(z) \tag{5b}
\end{align*}
$$

and the stationary value of $F$ is given by

$$
\begin{equation*}
F\left[\phi_{\mathrm{st}}^{\prime}, \phi_{v \mathrm{st}}\right]=\operatorname{Tr} \rho(z+L)^{-1} A_{v}=\left\langle A_{v}\right\rangle(z) \tag{6}
\end{equation*}
$$

Starting from the functional (4) and a suitable ansatz for $\phi^{\prime}$ and $\phi_{v}$ one can find approximations for $\left\langle A_{v}\right\rangle(z)$. This is in practice possible if one has some idea in which part of the Liouville space the stationary values (5) of $\phi^{\prime}$ and $\phi_{v}$ will be. Then this subset can be parametrized, and the parameters can be determined from $\delta F=0$.

In the next section we shall use linear parametrizations of the chosen subsets, in section 4 we will introduce nonlinear parametrizations as well. Common to all variational principles, the accuracy of the results will crucially depend on the chosen subsets. The variational method will be demonstrated in each section by well-known physical examples.

## 3. LINEAR VARIATIONAL ANSATZ

Let $\varphi_{v}$ and $\varphi_{v}^{\prime}$ be two sets of vectors such that linear combinations of $\varphi_{v}$ and of $\varphi_{v}^{\prime}$ lead to approximate expressions for $A_{v}(z)$ and $\rho(z)$, respectively. The question how to find these sets will be discussed later (see Section 3.2).

Now, according to the desired expression (5) for $\phi_{v}$ and $\phi^{\prime}$ it is reasonable to take the following ansatz:

$$
\begin{align*}
\phi_{v} & =\sum_{\mu} c_{v \mu} \varphi_{\mu}  \tag{7a}\\
\phi^{\prime} & =\sum_{\mu} c_{\mu}^{\prime} \varphi_{\mu}^{\prime} \tag{7b}
\end{align*}
$$

Inserting the ansatz (7) into (4) we find for $F$

$$
\begin{align*}
F\left[c^{\prime}, c\right]= & \sum_{\mu} c_{\mu}^{\prime} \operatorname{Tr} \varphi_{\mu}^{\prime} A_{v}+\sum_{\mu} c_{v \mu} \operatorname{Tr} \rho \varphi_{\mu} \\
& -\sum_{\mu, \lambda} c_{\mu}^{\prime} c_{v \lambda} \operatorname{Tr} \varphi_{\mu}^{\prime}(z+L) \varphi_{\lambda} \tag{8}
\end{align*}
$$

The stationary conditions of $F$ read

$$
\begin{align*}
\sum_{\lambda} c_{v \lambda} \operatorname{Tr} \varphi_{\mu}^{\prime}(z+L) \varphi_{\lambda}-\operatorname{Tr} \varphi_{\mu}^{\prime} A_{v} & =0  \tag{9a}\\
\sum_{\lambda} c_{\lambda}^{\prime} \operatorname{Tr} \varphi_{\lambda}^{\prime}(z+L) \varphi_{\mu}-\operatorname{Tr} \rho \varphi_{\mu} & =0 \tag{9b}
\end{align*}
$$

Equations (9) determine the stationary values

$$
\begin{gather*}
\left(c_{v \mu}\right)_{\mathrm{st}}=c_{v \mu}(z) \\
\left(c_{\mu}^{\prime}\right)_{\mathrm{st}}=c_{\mu}^{\prime}(z) \tag{10}
\end{gather*}
$$

which depend on the frequency $z$. Using the equations (9), we find for the stationary value of $F$

$$
\begin{equation*}
F_{\mathrm{st}}(z)=F\left[\phi_{\mathrm{st}}^{\prime}, \phi_{v \mathrm{st}}\right]=\operatorname{Tr} \rho \phi_{v \mathrm{st}}(z)=\operatorname{Tr} \phi_{\mathrm{st}}^{\prime}(z) A_{v} \tag{11}
\end{equation*}
$$

where

$$
\begin{align*}
\phi_{v \mathrm{st}}(z) & =\sum_{\mu}\left(c_{v \mu}\right)_{\mathrm{st}} \varphi_{\mu}  \tag{12}\\
\phi_{\mathrm{st}}^{\prime}(z) & =\sum_{\mu}\left(c_{\mu}^{\prime}\right)_{\mathrm{st}} \varphi_{\mu}^{\prime}
\end{align*}
$$

approximate $A_{v}(z)$ and $\rho(z)$, respectively. It suffices, however, to solve one of the decoupled systems of equations (9a) or (9b).

### 3.1. A Linear System of Differential Equations for Expectation Values

Transforming the equations (9a) or (9b) into the time region, and taking appropriate linear combinations, we find the following system of linear equations for $F_{\mathrm{st}}(t)=\left\langle A_{v}\right\rangle(t)$ :

$$
\begin{equation*}
\left\langle\dot{A}_{v}\right\rangle-i \sum_{\lambda}\left\langle A_{\lambda}\right\rangle \Lambda_{\lambda 0}=0 \tag{13a}
\end{equation*}
$$

where the kinetic matrix $\Lambda$ is defined by

$$
\begin{align*}
\Lambda & =\alpha^{-1} L \beta^{-1} \alpha \\
\alpha_{v \mu} & =\operatorname{Tr} \varphi_{v}^{\prime} A_{\mu}  \tag{13b}\\
\beta_{v \mu} & =\operatorname{Tr} \varphi_{v}^{\prime} \varphi_{\mu} \\
L_{v \mu} & =\operatorname{Tr} \varphi_{v}^{\prime} L \varphi_{\mu}
\end{align*}
$$

According to (9) the initial conditions are

$$
\begin{equation*}
\left\langle A_{v}\right\rangle(t=0)=\sum_{\mu} \operatorname{Tr} \rho \varphi_{\mu}\left(\beta^{-1} \alpha\right)_{\mu v} \tag{14}
\end{equation*}
$$

In (13) we have derived a system of equations for expectation values $\left\langle A_{v}\right\rangle$. Because of

$$
\begin{equation*}
\phi_{\mathrm{st}}^{\prime}(t)=\sum_{\mu}\left(c_{\mu}^{\prime}\right)_{\mathrm{st}}(t) \varphi_{\mu}^{\prime} \tag{15}
\end{equation*}
$$

we alternatively can use equations (9b) to derive a master equation for the relevant statistical operator $\rho_{\mathrm{rel}}(t)=\phi_{\mathrm{st}}^{\prime}(t)$. We will discuss this aspect later by means of a special example (see Section 3.6).

### 3.2. Conditions on $\varphi_{v}$ and $\varphi_{v}^{\prime}$

In equations (13) the kinetic matrix $\Lambda$ remains still undetermined as long as the vectors $\varphi_{v}$ and $\varphi_{v}^{\prime}$ are not specified. For a given set of observables $\left\{A_{v}\right\}$ and a given set of initial statistical operators $\{\rho\}$ we now discuss an appropriate choice of the subspaces $\left\{\varphi_{v}\right\}$ and $\left\{\varphi_{v}^{\prime}\right\}$.

It is physically reasonable to demand that the solution of (13) with initial conditions (14) would become exact for $t=0$. This implies

$$
\begin{equation*}
\left\langle A_{v}\right\rangle(t=0)=\operatorname{Tr} \rho A_{v} \tag{16}
\end{equation*}
$$

From (14) one sees that condition (16) is fulfilled if either

$$
\begin{equation*}
\rho=\sum_{\mu} r_{\mu} \varphi_{\mu}^{\prime} \tag{17}
\end{equation*}
$$

holds for all statistical operators $\rho$ to be considered or

$$
\begin{equation*}
A_{v}=\sum_{\mu} a_{v \mu} \varphi_{\mu} \tag{18}
\end{equation*}
$$

where $r_{\mu}$ and $a_{v \mu}$ are $c$ numbers.

### 3.3. Simplest Cholce of $\varphi_{v}$ and $\varphi_{v}^{\prime}$

The simplest ansatz, which fulfills the physically reasonable condition (16), is given by

$$
\begin{align*}
& \varphi_{v}=A_{v} \\
& \varphi_{v}^{\prime}=R_{v} \tag{19}
\end{align*}
$$

Here we have introduced a basis $\left\{R_{v}\right\}$ of the space spanned by the statistical operators $\rho$ which are to be considered:

$$
\begin{equation*}
\rho=\sum_{v} r_{v} R_{v} \tag{20}
\end{equation*}
$$

The ansatz (19) obviously fulfills (17) and (18) and, therefore, equation (16). This means: The solution (11) is exact for $z \rightarrow \infty$

$$
\begin{equation*}
F_{\mathrm{st}}(z \rightarrow \infty)=\left\langle A_{v}\right\rangle_{\mathrm{exact}}(z \rightarrow \infty) \tag{21}
\end{equation*}
$$

Because of $\beta=\alpha$ we find for the matrix $\Lambda$

$$
\begin{align*}
\Lambda & =\alpha^{-1} L \equiv \Omega \\
\alpha_{v \mu} & =\operatorname{Tr} R_{v} A \mu  \tag{22}\\
L_{v \mu} & =\operatorname{Tr} R_{v} L A_{\mu}
\end{align*}
$$

Inspection of the equation of motion (13a) shows that the result (22) for the matrix $\Lambda$ implies the correct values for the derivatives $\left\langle\dot{A}_{v}\right\rangle$ at $t=0$. The solution of (13a), therefore, represents an extrapolation of the short-time behavior of $\left\langle A_{v}\right\rangle(t)$. Normally, this procedure does not lead to a description of damping phenomena. To describe such effects it is necessary to find a solution of $\left\langle A_{v}\right\rangle(t)$ which is meaningful for long times, too. We will treat this problem in the next section.

### 3.4. Cholce of $\varphi_{v}$ and $\varphi_{v}^{\prime}$ for "Slow Varlables"

Now let us seek a more useful ansatz, which would give the exact result for two different values of $z$, say $z \rightarrow \infty$ and $z=z_{0}$. We already know that the ansatz $\varphi_{v}^{\prime}=R_{v}$ gives the exact result for $\left\langle A_{v}\right\rangle(z)$ at $z \rightarrow \infty$, $F_{\mathrm{st}}(z \rightarrow \infty)=\left\langle A_{v}\right\rangle_{\text {exact }}(z \rightarrow \infty)$, whatever the ansatz for $\varphi_{v}$ may be. Vice versa, the ansatz $\varphi_{v}=A_{v}\left(z_{v}\right)$ yields the exact result for $\left\langle A_{v}\right\rangle(z)$ at $z=z_{v}$, $F_{\mathrm{st}}\left(z_{v}\right)=\left\langle A_{v}\right\rangle_{\mathrm{exact}}\left(z_{v}\right)$, whatever we may have used for $\varphi_{v}^{\prime}$. This follows from (9b) and (11). We therefore expect that the ansatz

$$
\begin{align*}
& \varphi_{v}^{\prime}=R_{v} \\
& \varphi_{v}=A_{v}\left(z_{v}\right) \tag{23}
\end{align*}
$$

would give an appropriate interpolation for $\left\langle A_{v}\right\rangle(z)$ between the exact values $\left\langle A_{v}\right\rangle\left(z_{v}\right)$ and $\left\langle A_{v}\right\rangle(z \rightarrow \infty)$.

In order for this interpolation between the two frequencies $z=\infty$ and $z=z_{v}$ to represent a useful approximation of the exact result, we have to postulate some properties of the observables $A_{v}$ : Let us assume that $e^{-i \omega_{v} t} A_{v}(t)$ is a "slowly varying quantity," which means that its Fourier transform is peaked at $\omega=0$, or the Fourier transform of $A_{v}(t)$ is peaked at $\omega=-\omega_{v}$, respectively. If $e^{-i \omega_{v} t} A_{v}(t)$ were constant in time, we would have

$$
\begin{equation*}
A_{v}(z)=c_{v}(z) A_{v}\left(z_{v}\right) \tag{24}
\end{equation*}
$$

so that the ansatz

$$
\begin{equation*}
\varphi_{v}=A_{v}\left(z_{v}\right) \tag{25}
\end{equation*}
$$

would lead to the exact result for $\left\langle A_{v}\right\rangle(z)$. Hence we expect for "slowly varying" $e^{-i \omega_{v} t} A_{v}(t)$ the variational ansatz $\varphi_{v}=A_{v}\left(z_{v}\right)$ to be adequate.

To find the frequency $\omega_{v}$ in the Fourier spectrum of $A_{v}(t)$ we use perturbation theory and demand that there exists a main part $L_{0}$ of the Liouvillian $L=L_{0}+g L_{1}$ with

$$
\begin{equation*}
L_{0} A_{v}=\omega_{v} A_{v} \tag{26}
\end{equation*}
$$

In this case we choose $z_{v}=-\omega_{v}+i \eta$. Then equation (26) leads to

$$
\begin{equation*}
A_{v}\left(z_{v}\right)=-i \int_{0}^{\infty} d \tau e^{i z_{v} \tau} e^{i L \tau} A_{v}=-i \int_{0}^{\infty} d \tau e^{-\eta \tau} e^{i L \tau} e^{-i L_{0} \tau} A_{v} \tag{27}
\end{equation*}
$$

Since we will use the linear ansatz

$$
\begin{equation*}
\phi_{v}=\sum_{\mu} c_{v \mu} \varphi_{\mu}=\sum_{\mu} c_{v \mu} A_{\mu}\left(z_{\mu}\right) \tag{28}
\end{equation*}
$$

we can take linear combinations of $A_{v}$ in (27) as well. From this we conclude the following:

Given a set $\left\{A_{v}\right\}$ invariant under $L_{0}$

$$
\begin{equation*}
L_{0} A_{v}=\sum_{\mu} L_{v \mu}^{0} A_{\mu} \tag{29}
\end{equation*}
$$

then a suitable ansatz for $\varphi_{v}^{\prime}$ and $\varphi_{v}$ is

$$
\begin{align*}
& \varphi_{v}^{\prime}=R_{v} \\
& \varphi_{v}=-i \int_{0}^{\infty} d \tau e^{-\eta \tau} e^{i L \tau} e^{-i L_{0} \tau} A_{v} \tag{30}
\end{align*}
$$

Now ansatz (30) leads to the following result for the matrix $\Lambda$ defined in (13b):

$$
\begin{equation*}
\Lambda=\Omega+i \Gamma \tag{31}
\end{equation*}
$$

where the expression for the matrix $\Omega$ reads

$$
\begin{align*}
(\alpha \Omega)_{v \mu} & =\operatorname{Tr} R_{v} L A_{\mu} \\
\alpha_{v \mu} & =\operatorname{Tr} R_{v} A_{\mu} \tag{32}
\end{align*}
$$

This term has already been found by means of the simple ansatz (19) [compare Eq. (22)], where it was the only contribution to $\Lambda$. The matrix $\Gamma$ is expressed as

$$
\begin{equation*}
(\alpha \Gamma)_{v \mu}=g^{2} \int_{0}^{\infty} d \tau e^{-\eta \tau} \operatorname{Tr} R_{v} L_{1}^{\prime} e^{i L_{0} \tau} L_{1} e^{-i L_{0} \tau} A_{\mu}+O\left(g^{3}\right) \tag{33}
\end{equation*}
$$

where $\alpha$ is given by (32) and $L_{1}^{\prime}$ is defined by

$$
\begin{equation*}
L_{1}^{\prime} X=L_{1} X-\sum_{\mu, v} L_{1} A_{\mu}\left(\alpha^{-1}\right)_{\mu \nu} \operatorname{Tr} R_{v} X \tag{34}
\end{equation*}
$$

To arrive at formula (33) we have used the result of first-order perturbation theory for the time-evolution operator $e^{i L t}$. ${ }^{2}$

The basic results of our linear variational method are presented in equations (9) and (11). These equations can be rewritten as a linear system of differential equations (13) for the expectation values of observables $A_{v}$. The kinetic matrices involved may be specialized to yield the expressions (22) or (31), (32) and (33) respectively. These variational results comprise various projection-operator approaches. This fact will be discussed in the Appendix. Here we only want to illustrate this connection by simple examples.

### 3.5. Example: Linear Relaxation of Parallel Magnetization

First let us discuss the relaxation of a magnetization $M_{z}$ parallel to a static magnetic field $H$ if the field is changed by a small amount $\Delta H$. The purpose of this example is to illustrate the connection to a Markovian approximation of Mori's theory ${ }^{(4)}$ of linear Langevin equations. Consider a system described by the Hamiltonian

$$
\begin{align*}
\mathscr{H} & =\mathscr{K}_{0}+g \mathscr{K}_{1} \\
\mathscr{K}_{0} & =-H M_{z}+\mathscr{K}_{0}(H=0), \quad\left[M_{z}, \mathscr{K}_{0}(0)\right]=0 \tag{35}
\end{align*}
$$

and the statistical operator at $t=0$

$$
\begin{align*}
\rho & =R(\Delta H)=[Z(\Delta H)]^{-1} e^{-\beta\left(\Upsilon \mathcal{X}-\Delta H M_{z}\right)}, \\
& =R(0)+\left.\Delta H \frac{\partial R}{\partial \Delta H}\right|_{\Delta H=0} \tag{36}
\end{align*}
$$

for small $\Delta H$. Thus the space of considered statistical operators $\{\rho\}$ is spanned by $R_{1}=R(0)$ and $R_{2}=\left.(\partial R / \partial \Delta H)\right|_{\Delta H=0}$. We therefore choose

$$
\begin{align*}
& A_{1}=1, \quad A_{2}=M_{z}-\left\langle M_{z}\right\rangle_{\infty}, \quad\left\langle M_{z}\right\rangle_{\infty}=\operatorname{Tr} R(0) M_{z} \\
& R_{1}=R(0), \quad R_{2}=\left.\frac{\partial R}{\partial \Delta H}\right|_{\Delta H=0} \tag{37}
\end{align*}
$$

Since equation (29) is fulfilled, we can consider $A_{1}$ and $A_{2}$ to be "slow variables." Applying ansatz (30) we immediately find from (32), (33), and (13a): $\Omega=0, \alpha_{11}=1, \alpha_{22}=\chi$, and

$$
\begin{align*}
\left\langle\dot{M}_{z}\right\rangle & =-\Gamma\left(\left\langle M_{z}\right\rangle-\left\langle M_{z}\right\rangle_{\infty}\right) \\
\left\langle M_{z}\right\rangle(0) & =\left\langle M_{z}\right\rangle_{\infty}+\Delta H_{\chi} \tag{38}
\end{align*}
$$

[^1]where the damping constant $\Gamma$ is given by
\[

$$
\begin{align*}
\Gamma & =\left.\chi^{-1}\left(\frac{g}{\hbar}\right)^{2} \int_{0}^{\infty} d \tau e^{-\eta \tau} \operatorname{Tr} \frac{\partial R}{\partial \Delta H}\right|_{\Delta H=0}\left[\mathscr{H}_{1},\left[\mathscr{K}_{1}(\tau), M_{z}\right]\right]  \tag{39}\\
\mathscr{K}_{1}(\tau) & =e^{(i / \hbar) \mathscr{K}_{0} \tau \mathcal{H}_{1} e^{-(i / \hbar) \mathscr{K}_{0} \tau}}
\end{align*}
$$
\]

For convenience we have introduced the isothermal magnetic susceptibility $\chi=\left(\partial M_{z} / \partial H\right)_{T}$. Equations (38) and (39) are usually derived by Mori's theory ${ }^{(4)}$ and are a standard result in linear magnetic relaxation theory.

### 3.6. Example: Master Equation for the Statistical Operator $R_{s}$ of a Subsystem

In our second example we consider a system $\mathscr{K}_{s}$ in contact with a heat reservoir $\mathscr{H}_{B}$. Let the Hamiltonian be

$$
\begin{equation*}
\mathfrak{H}=\mathscr{K}_{s}+\mathscr{K}_{B}+g \mathscr{K}_{1} \tag{40}
\end{equation*}
$$

The equilibrium statistical operator of the heat reservoir is denoted by

$$
\begin{equation*}
R_{B}=Z_{B}^{-1} e^{-\beta \mathscr{K}_{B}} \tag{41}
\end{equation*}
$$

We now use ansatz (30) and equation (9b) to derive the standard master equation for the relevant statistical operator $R_{s}$ of the subsystem $\mathscr{H}_{s}$. Let

$$
\begin{align*}
& A_{v}=A_{v}^{(s)} \\
& R_{v}=R_{B} \cdot R_{v}^{(s)} \tag{42}
\end{align*}
$$

where both sets $\left\{A_{v}^{(s)}\right\}$ and $\left\{\boldsymbol{R}_{v}^{(s)}\right\}$ span a complete basis of the Liouville space of the subsystem. Then condition (29) is fulfilled and ansatz (30) immediately yields the following differential equation for $R_{s}{ }^{3}$ :

$$
\begin{align*}
R_{s}(t) & =\operatorname{Tr}_{B} \phi^{\prime}(t)=\sum_{v} c_{v}^{\prime}(t) R_{v}^{(s)} \\
\dot{R}_{s} & =-i L_{s} R_{s}-\left(\frac{g}{\hbar}\right)^{2} \int_{0}^{\infty} d \tau e^{-\eta \tau} \operatorname{Tr}_{B}\left[\mathscr{K}_{1},\left[\mathscr{K}_{1}(-\tau), R_{B} R_{s}\right]\right] \tag{43}
\end{align*}
$$

which is the standard result for the time evolution of the reduced density operator. ${ }^{(5)}$

## 4. NONLINEAR VARIATIONAL ANSATZ

In Section 3 we have considered linear parametrizations of the various sets $\left\{\phi^{\prime}\right\}$ and $\left\{\phi_{v}\right\}$. They lead to linear systems of equations for expectation values of macroscopic observables $A_{v}$, or to linear master equations for relevant statistical operators $R_{s}$, respectively. In this section we want to

[^2]discuss a nonlinear parametrization of the set $\left\{\phi^{\prime}\right\}$. We, therefore, drop the condition that $\phi^{\prime}$ be only linearly dependent on the variational parameters $\left\{c_{v}^{\prime}\right\}$ and write
\[

$$
\begin{align*}
\phi^{\prime}(z) & =-i \int_{0}^{\infty} d t e^{i z t} \phi^{\prime}\left(c_{v}^{\prime}(t)\right) \\
\phi_{v} & =\sum_{\mu} c_{v \mu} \varphi_{\mu} \tag{44}
\end{align*}
$$
\]

The nonlinear function $\phi^{\prime}\left(c_{v}^{\prime}\right)$ will be defined later. Inserting the ansatz (44) into (4) we find for our functional

$$
\begin{equation*}
F\left[c^{\prime}, c\right]=\operatorname{Tr} \phi^{\prime}(z) A_{v}+\sum_{\mu} c_{v \mu} \operatorname{Tr} \rho \varphi_{\mu}-\sum_{\mu} c_{v \mu} \operatorname{Tr} \phi^{\prime}(z)(z+L) \varphi_{\mu} \tag{45}
\end{equation*}
$$

Variation with respect to the parameters $c_{c \mu}$ gives

$$
\begin{equation*}
\operatorname{Tr} \phi^{\prime}(z)(z+L) \varphi_{\mu}-\operatorname{Tr} \rho \varphi_{\mu}=0 \tag{46}
\end{equation*}
$$

These equations determine the stationary values of $c_{v}^{\prime}$, which will be denoted by $\lambda_{\mathrm{r}}(t)$

$$
\begin{equation*}
\left(c_{v}^{\prime}(t)\right)_{\mathrm{st}}=\lambda_{v}(t) \tag{47}
\end{equation*}
$$

Using relation (46) and the definition

$$
\begin{equation*}
R(z)=\left(\phi^{\prime}(z)\right)_{\mathrm{st}} \tag{48}
\end{equation*}
$$

we can write for the stationary value of $F$

$$
\begin{equation*}
F_{\mathrm{st}}=\operatorname{Tr} R(z) A_{v} \tag{49}
\end{equation*}
$$

Transforming equations (46) into the time region, we find for $R(t)=$ $\phi^{\prime}\left(\lambda_{v}(t)\right)$ the following differential equations:

$$
\begin{equation*}
\operatorname{Tr} \varphi_{\mu}(\dot{R}+i L R)=0 \tag{50}
\end{equation*}
$$

which have to be solved with the initial conditions

$$
\begin{equation*}
\operatorname{Tr} \varphi_{\mu} R(t=0)=\operatorname{Tr} \varphi_{\mu} \rho \tag{51}
\end{equation*}
$$

Equations (50) and (51) are equivalent to (46).
We now want to choose an explicit form of the function $\phi^{\prime}\left(c_{v}^{\prime}\right)$. For a given set of macroscopic observables $A_{v}$ we will consider the following set of initial statistical operators:

$$
\begin{equation*}
\rho\left(\lambda_{v}^{0}\right)=\exp \left(-\sum_{v} \lambda_{v}^{0} A_{v}\right) \tag{52}
\end{equation*}
$$

where the parameters $\lambda_{v}^{0}$ are considered to be given by the expectation values of $A_{v}$ at $t=0$.

In the following we will suppose that the unity operator is a linear combination of $A_{v}$ :

$$
\begin{equation*}
1=\sum_{v} a_{v} A_{v} \tag{53}
\end{equation*}
$$

From the form (52) we are led to take the following function $\phi^{\prime}\left(c_{v}^{\prime}\right)$ as an ansatz in (44):

$$
\begin{equation*}
\phi^{\prime}\left(c_{v}^{\prime}\right)=\exp \left(-\sum_{v} c_{v}^{\prime} A_{v}\right) \tag{54}
\end{equation*}
$$

The "relevant statistical operator" (49) then reads

$$
\begin{equation*}
R(t)=\exp \left[-\sum_{v} \lambda_{v}(t) A_{v}\right] \tag{55}
\end{equation*}
$$

The parameters $\lambda_{v}(t)$ are the stationary values of $c_{v}^{\prime}(t)$, and are to be calculated from equations (50) together with the initial condition (51): $\lambda_{v}(0)=\lambda_{v}^{0}$.

Let us summarize: Our nonlinear variational ansatz (44) and (54) has led to the system of equations (50) for $R(t)$, the form of which is given by (55). Choosing the same set $\varphi_{v}$ as in Section 3, we now can derive a closed system of nonlinear equations for the expectation values $\left\langle A_{v}\right\rangle(t)$. This procedure will be outlined in the following section.

### 4.1. A Nonlinear System of Differential Equations for Expectation Values

Equation (49) defines $\left\langle A_{v}\right\rangle(t)$ as a function of $\lambda_{v}(t)$ :

$$
\begin{equation*}
\left\langle A_{v}\right\rangle(t)=\left\langle A_{v}\right\rangle\left(\lambda_{\mu}(t)\right) \tag{56}
\end{equation*}
$$

Conversely, the stationary functions $\lambda_{v}(t)$ and the statistical operator $R(t)$ (55) depend on the time $t$ by the functions $\left\langle A_{v}\right\rangle(t): R(t)=R\left(\left\langle A_{v}\right\rangle(t)\right)$. Now because of relation (53) $R(t)$ may be expressed as ${ }^{(6)}$

$$
\begin{equation*}
R=\sum_{\lambda}\left\langle A_{\lambda}\right\rangle \frac{\partial R}{\partial\left\langle A_{\lambda}\right\rangle} \tag{57}
\end{equation*}
$$

By means of the identity (57) we can transform equations (50) into the following nonlinear differential equations:

$$
\begin{align*}
&\left\langle\dot{A}_{v}\right\rangle-i \sum_{\lambda}\left\langle A_{\lambda}\right\rangle \Lambda_{\lambda v}\left(\left\langle A_{\kappa}\right\rangle\right)=0^{4} \\
& \Lambda=L \beta^{-1} \\
& L_{v \mu}=\operatorname{Tr} \frac{\partial R}{\partial\left\langle A_{v}\right\rangle} L \varphi_{\mu}  \tag{58}\\
& \beta_{v \mu}=\operatorname{Tr} \frac{\partial R}{\partial\left\langle A_{v}\right\rangle} \varphi_{\mu}
\end{align*}
$$

${ }^{4}$ An equivalent form to the equation of motion (58) reads

$$
\left\langle\dot{A}_{v}\right\rangle-i \sum_{\lambda} \operatorname{Tr} R L \varphi_{\lambda}\left(\beta^{-1}\right)_{\lambda v}=0
$$

The system (58) is the analog to the system (13). Contrary to equations (13), however, (58) generally constitutes a system of nonlinear differential equations, because the kinetic matrix $\Lambda$ depends on $\left\langle A_{v}\right\rangle$.

It holds that $R(t=0)=\rho$, hence it follows from the relation $F_{\text {st }}=$ $\operatorname{Tr} R(t) A_{v}$, that $F_{\mathrm{st}}(t)$ becomes the exact value of $\left\langle A_{v}\right\rangle(t)=0$, whatever the special choice for $\varphi_{v}$ may be.

In section 3, possible choices of $\varphi_{v}$ were discussed. Let us use here the same two sets

$$
\begin{equation*}
\varphi_{v}=A_{v} \tag{59a}
\end{equation*}
$$

and

$$
\begin{equation*}
\varphi_{v}=-i \int_{0}^{\infty} d \tau e^{-\eta \tau} e^{i L \tau} e^{-i L_{0} \tau} A_{v} \tag{59b}
\end{equation*}
$$

If it holds that $L_{0} A_{v}=\omega_{v} A_{v}$, then the $\varphi_{v}$ of (59b) again lead to the exact values of $\left\langle A_{v}\right\rangle(z)$ at $z=z_{v}=-\omega_{v}+i \eta$. This fact follows from (46), (48), and (49):

$$
\begin{equation*}
\operatorname{Tr} \rho A_{v}\left(z_{v}\right)=\operatorname{Tr} R\left(z_{v}\right) A_{v}=F_{\mathrm{st}}\left(z_{v}\right) \tag{60}
\end{equation*}
$$

Applying the $\varphi_{v}$ of (59) we can rewrite the kinetic matrix (58). In case (59a) we find

$$
\begin{align*}
\beta & =1 \\
\Lambda & =\Omega  \tag{61}\\
\Omega_{v \mu} & =\operatorname{Tr} \frac{\partial R}{\partial\left\langle A_{v}\right\rangle} L A_{\mu}
\end{align*}
$$

In case (59b) we have

$$
\begin{equation*}
\Lambda=\Omega+i \Gamma \tag{62}
\end{equation*}
$$

where $\Omega$ has already been defined by (61) and $\Gamma$ is given by

$$
\begin{align*}
\Gamma_{v \mu} & =g^{2} \int_{0}^{\infty} d \tau e^{-\eta \tau} \operatorname{Tr} \frac{\partial R}{\partial\left\langle A_{v}\right\rangle} L_{1}^{\prime}(t) e^{i L_{0} \tau} L_{1} e^{-i L_{0} \tau} A_{\mu}+0\left(g^{3}\right) \\
L_{1}^{\prime}(t) X & =L_{1} X-\sum_{\mu} L_{1} A_{\mu} \operatorname{Tr} \frac{\partial R}{\partial\left\langle A_{\mu}\right\rangle} X \tag{63}
\end{align*}
$$

To illustrate the nonlinear formalism outlined above, we will apply it to some simple examples in the next sections.

### 4.2. Example: Time-Dependent Hartree-Fock Equation

As a first example for the nonlinear system (58) let us apply ansatz (59a). Consider a quantum gas of bosons or fermions with two-particle
interactions:

$$
\begin{align*}
\mathscr{H}_{0} & =\sum_{1} \epsilon_{1} a_{1}^{+} a_{1} \\
\mathscr{H}_{1} & =\frac{1}{2 N} \sum_{11^{\prime} 2^{\prime}} V_{121^{\prime} 2^{\prime}} a_{1}^{+} a_{2}^{+} a_{2^{\prime}} a_{1^{\prime}} \tag{64}
\end{align*}
$$

We choose

$$
\begin{align*}
A_{0} & =1 \\
A_{v} & =a_{1}^{+} a_{2}  \tag{65}\\
\phi^{\prime}\left(c_{0}^{\prime}, c_{12}^{\prime}\right) & =\exp \left(-\sum_{12} c_{12}^{\prime} a_{1}^{+} a_{2}-c_{0}^{\prime}\right)
\end{align*}
$$

The kinetic matrix (61) can be evaluated explicitly. From (58) we arrive at the following nonlinear system of equations for the one-particle density matrix $\left\langle a_{1}^{+} a_{2}\right\rangle(t)=\langle 2| \rho^{(1)}(t)|1\rangle$

$$
\begin{align*}
& \frac{\partial}{\partial t}\langle 2| \rho^{(1)}(t)|1\rangle=-(i / \hbar)\left(\epsilon_{2}-\epsilon_{1}\right)\langle 2| \rho^{(1)}(t)|1\rangle \\
&+(i / \hbar) \frac{1}{2 N} \sum_{344^{\prime}}\left\{\left(V_{424^{\prime} 3} \pm V_{244^{\prime} 3}\right)\left\langle 4^{\prime}\right| \rho^{(1)}(t)|4\rangle\langle 3| \rho^{(1)}(t)|1\rangle\right. \\
&-\left(V_{434^{\prime} 1} \pm V_{344^{\prime} 1}\right)\left\langle 4^{\prime}\right| \rho^{(1)}(t)|4\rangle \\
&\left.\times\langle 2| \rho^{(1)}(t)|3\rangle\right\} \tag{66}
\end{align*}
$$

They are the well-known Hartree-Fock equations. ${ }^{(7)}$

### 4.3. Example: Nonlinear Heat Conduction

To discuss ansatz (59b) let us consider the heat conduction between two reservoirs $\mathscr{K}_{1}$ and $\mathscr{K}_{2}$

$$
\begin{equation*}
\mathscr{H}=\mathscr{H}_{1}+\mathscr{K}_{2}+g \mathscr{H}_{12}, \quad\left[\mathscr{K}_{1}, \mathscr{H}_{2}\right]=0 \tag{67}
\end{equation*}
$$

We take

$$
\begin{align*}
A_{0} & =1 \\
A_{v} & =\mathscr{K}_{v}, \quad v=1,2  \tag{68}\\
\phi^{\prime}\left(c_{v}^{\prime}\right) & =\exp \left(-\sum_{v} c_{v}^{\prime} \mathscr{K}_{v}-c_{0}^{\prime}\right)
\end{align*}
$$

From (50) and (59b) we directly find the equation for heat conduction

$$
\begin{equation*}
\left\langle\dot{\mathscr{K}}_{1}\right\rangle=-\Gamma\left(\beta_{v}\right)\left(\beta_{2}-\beta_{1}\right)=-\left\langle\dot{\mathscr{C}}_{2}\right\rangle \tag{69}
\end{equation*}
$$

where the thermal conductivity $\Gamma\left(\beta_{v}\right)$ is given by

$$
\begin{align*}
& \Gamma\left(\beta_{v}\right)=\int_{0}^{\infty} d \tau e^{-m} \int_{0}^{1} d \lambda \operatorname{Tr} R^{1-\lambda \dot{\mathcal{C}}} \dot{\mathcal{C}}_{1}(\tau) R^{\lambda \dot{\mathcal{C}}} \dot{\mathcal{C}}_{1}+O\left(g^{3}\right) \\
& R\left(\beta_{v}\right)=\left[Z\left(\beta_{v}(t)\right)\right]^{-1} \exp \left(-\sum_{v} \beta_{v}(t) \mathscr{K}_{v}\right)  \tag{70}\\
& \dot{K}_{1}=(i / \hbar) g\left[\mathscr{K}_{12}, \mathscr{K}_{1}\right] \\
& \mathscr{H}(\tau)=e^{(i / \hbar)\left(\mathscr{X}_{1}+\mathscr{X}_{2}\right) \tau \dot{\mathcal{C}}_{1} e^{(-i / \hbar)\left(\mathscr{X}_{1}+\mathscr{X}_{2}\right) \tau}}
\end{align*}
$$

The thermal conductivity $\Gamma$ itself depends on $\beta_{v}(t)$; thus (69) represents a nonlinear system of heat conduction equations. ${ }^{(8)}$

### 4.4. Example: Nonlinear Dynamics of a Spin System

In Sections 4.2 and 4.3 we have regarded examples where the kinetic matrices $\Lambda$ were given by $\Lambda=\Omega$ or by $\Lambda=i \Gamma$, respectively. We now want to treat a simple example where both terms $\Omega$ and $\Gamma$ contribute to $\Lambda$.

Consider the physical system of Section 3.6, where we now specialize the subsystem $S$ to be constituted of $N$ spins $s=1 / 2$ in an external magnetic field $H$. Thus we write

$$
\begin{align*}
\mathscr{K}_{s} & =-\gamma H \sum_{i} S_{i}^{z} \\
g \mathscr{K}_{1} & =g \sum_{i}\left(B_{i}^{+} S_{i}^{-}+B_{i}^{-} S_{i}^{+}\right)  \tag{71}\\
\left\langle B_{i}^{ \pm}\right\rangle & =\operatorname{Tr}_{B} R_{B} B_{i}^{ \pm}=0, \quad R_{B}=Z_{B}^{-1} e^{-\beta \mathscr{K}_{B}}
\end{align*}
$$

Now let us take

$$
\begin{align*}
A_{0} & =1 \\
A_{v} & =\sum_{i} S_{i}^{v}, \quad v=+,-, z  \tag{72}\\
\phi^{\prime}\left(c_{0}^{\prime}, c_{v}^{\prime}\right) & =R_{B} \cdot \exp \left(-\sum_{v} c_{v}^{\prime} \sum_{i} S_{i}^{v}-c_{0}^{\prime}\right)
\end{align*}
$$

and use ansatz (59b). Then a straightforward calculation of (58) yields the following nonlinear equation for $\mathbf{M}=\gamma \sum_{i}\left\langle\mathbf{S}_{i}\right\rangle$ :

$$
\begin{equation*}
\dot{\mathbf{M}}=-\gamma(\mathbf{M} \times \mathbf{H})-D \cdot\left(\mathbf{M}-\mathbf{M}_{\infty}\right)-\gamma^{\prime} \mathbf{M} \times M_{z} \mathbf{e}_{z}+\gamma^{\prime \prime}(\mathbf{M} \times(\mathbf{M} \times \mathbf{H})) \tag{73}
\end{equation*}
$$

where

$$
\left.\begin{array}{c}
\mathbf{H}=\left[\begin{array}{l}
0 \\
0 \\
H
\end{array}\right] \mathbf{M}_{\infty}=\left[\begin{array}{c}
0 \\
0 \\
M_{\infty}
\end{array}\right], \quad M_{\infty}=\frac{N \gamma \hbar}{2} \tanh \frac{\gamma \hbar H}{2 k T} \\
D \\
0
\end{array} \quad 1 / \tau \quad \begin{array}{ll}
1 / \tau & 0 \\
0 & 2 / \tau
\end{array}\right] \quad \begin{gathered}
\gamma \gamma^{\prime}=-2 g^{2} \operatorname{Re}\left(\frac{1}{N^{2}} \sum_{i \neq j} \chi_{i j}(\gamma H)\right)  \tag{73b}\\
\gamma H \gamma^{\prime \prime}=-2 g^{2} \operatorname{Im}\left(\frac{1}{N^{2}} \sum_{i \neq j} \chi_{i j}(\gamma H)\right) \\
\gamma \tau^{-1}=-\left(g^{2} / 2\right) M_{\infty}^{-1} \operatorname{Im}\left(\sum_{i} \chi_{i i}(\gamma H)\right) \\
\chi_{i j}(\omega)=\frac{1}{\hbar} \int_{0}^{\infty} d t e^{i \omega t}\left\langle\left[B_{i}^{+}(t), B_{j}^{-}\right]\right\rangle
\end{gathered}
$$

The first term on the right of equation (73) constitutes the linear reversible motion, the second term is the Bloch damping, the third renormalizes the reversible motion by a molecular field:

$$
\begin{align*}
(\dot{\mathbf{M}})_{\mathrm{rev}} & =-\gamma\left(\mathbf{M} \times \mathbf{H}_{\mathrm{eff}}\right) \\
\mathbf{H}_{\mathrm{eff}} & =\left(H+\frac{\gamma^{\prime}}{\gamma} M_{z}\right) \mathbf{e}_{z} \tag{74}
\end{align*}
$$

and the last term is the Landau-Lifschitz damping ${ }^{(9)}$.

## 5. CONCLUSIONS

It has been shown that the Schwinger variational principle of scattering theory is suited for quantum statistics, where it allows a derivation of macroscopic equations of motion.

The variational method might be useful to improve known results by optimizing parameters, or will make it possible to attack complicated problems if the dynamics for different limiting cases is known and one wants to treat an intermediate situation.

## APPENDIX

To point out the connection of our result of linear relaxation problems to standard "projection-operator" approaches, let us define the following
idempotent operator $P$ :

$$
\begin{equation*}
P X=\sum_{v, \mu} A_{v}\left(\alpha^{-1}\right)_{v \mu} \operatorname{Tr} R_{\mu} X \tag{A.1}
\end{equation*}
$$

where $\alpha$ is given by (22). Then by means of (A.1), relation (11) can be cast into the following form:
(a) In case of ansatz (19) one obtains

$$
\begin{equation*}
F_{\mathrm{st}}(z)=\operatorname{Tr} \rho(z+P L P)^{-1} A_{v} \tag{A.2}
\end{equation*}
$$

(b) In case of ansatz (30) and $\mathrm{L}_{0} \mathrm{~A}_{v}=0$ for all $v$ one arrives at

$$
\begin{equation*}
F_{\mathrm{st}}(z)=\operatorname{Tr} \rho[z+P L P-M(i \eta)]^{-1} A_{v} \tag{A.3}
\end{equation*}
$$

where

$$
\begin{equation*}
M(z)=P L Q(z+Q L Q)^{-1} Q L P, \quad Q=1-P \tag{A.4}
\end{equation*}
$$

Inspection of the results (A.2) and (A.3), (A.4) and the exact expression for $\left\langle A_{v}\right\rangle(z)$

$$
\begin{equation*}
\left\langle A_{v}\right\rangle(z)=\operatorname{Tr} \rho(z+P L P-M(z))^{-1} A_{v} \tag{A.5}
\end{equation*}
$$

shows that the ansatz (19) in the variational method neglects memory effects, but gives an exact description of the systematic part of dynamics, whereas ansatz (30) corresponds to a Markovian approximation in the memory kernels.

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[^0]:    ${ }^{1}$ Sonderforschungsbereich 65 Darmstadt-Frankfurt, Frankfurt, West Germany.

[^1]:    ${ }^{2}$ In the case where $L_{1}^{\prime}=L_{1}$ and equation (26) holds, we directly find expression (33) for $\Gamma$ by the ansatz

    $$
    \varphi_{v}=\left\{\left(z_{v}+L_{0}\right)^{-3}-\left(z_{v}+L_{0}\right)^{-1} L_{1}\left(z_{v}+L_{0}\right)^{-1}\right\} A_{v}
    $$

[^2]:    ${ }^{3} \mathrm{We}$ have taken $\left\langle\mathcal{K}_{1}\right\rangle_{B}=\mathrm{Tr}_{B} R_{B} \mathscr{K}_{1}=0$.

