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# Dynamics of the Dicke laser model II: an example of a semi-classical theory in the presence of focal points 

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#### Abstract

The nature of the time-dependent WKB-Maslov approximation is investigated for the quantum mechanical system which is equivalent to the Dicke Laser model. The local semi-classical solution and its continuation beyond the focal points are discussed.


## 1. Introduction

The previous paper [1] on the dynamics of the Dicke model was concerned with explicit calculations of the most important dynamical quantities of the system. The main tool in these calculations was the time-dependent WKB-Maslov method. In the present paper this method will be studied in detail; in particular, the peculiarities that appear in the application to the Dicke Hamiltonian will be analysed. This is a subject of its own interest because it gives an interesting example for the use of the WKB method. This example shows how miserable our imagination is, even for quantum systems with one degree of freedom, if trained only with the harmonic oscillator.

The quantum mechanical problem to be discussed is defined in the next section. Its relation to the Dicke Hamiltonian was discussed before [1] and is unimportant in the present context, so that the subject is fully understandable without any reference to quantum optics. The corresponding classical mechanical problem is then considered. The classical solution shows clearly the local nature of the Hamilton-Jacobi theory: The Cauchy problem for the Hamilton-Jacobi equation has a unique solution only for small times $t<T_{0}$. For larger times the solution becomes multivalued, the branch-points are the focal points. As a consequence, the semi-classical WKB solution can be derived only for $t<T_{0}$ (Section 3), and in fact, it diverges at the focal point $T_{0}$. The main problem is the continuation of this local semi-classical solution beyond the focal points. This can be done according to a method of Maslov [2] by going over to the momentum space. Then, beyond a focal point new branches $S_{k}(x, t)$ of the action contribute to the semi-classical solution. In this way, the global semi-classical solution is obtained (Section 4). It is valid on a much larger time interval, except in the neighborhood of the focal points. Since the number of branches $S_{k}(x, t)$ increases linearly with time $t$, the error in the semi-classical solution increases in the same way. This is the mechanism which finally leads to the breakdown of the global semi-classical solution for large times, where the system then shows an essentially quantum-mechanical behavior.

## 2. The quantum mechanical and the classical Cauchy problems

We will consider the following time-dependent Schrödinger equation

$$
\begin{align*}
i h \frac{\partial y}{\partial t} & =-\frac{h^{2}}{2}\left(x \frac{\partial^{2} y}{\partial x^{2}}+\frac{\partial^{2}}{\partial x^{2}} x y\right)+V(x) y  \tag{2.1}\\
& \xlongequal{\text { def }} H y
\end{align*}
$$

with the potential energy

$$
\begin{equation*}
V(x)=\frac{1}{4} x^{3}-\eta x+\frac{\lambda^{2}}{x}, \quad \eta, \lambda>0 \tag{2.2}
\end{equation*}
$$

For simplicity, a (non-resonant) term $\delta x^{2} / 2$ in $V(x)$ [1] has been dropped. The whole analysis goes through with some minor changes in this non-resonant case too.

The Hamiltonian

$$
\begin{equation*}
H=-h^{2} \frac{\partial}{\partial x} x \frac{\partial}{\partial x}+V(x) \tag{2.3}
\end{equation*}
$$

defines a symmetric differential operator $H$ in $L^{2}(0, \infty)$. It is essential for the formulation of a Cauchy problem for equation (2.1) to extend $H$ to a self-adjoint operator. Both end-points $x=0$ and $\infty$ are singular points which belong, according to the classification of H . Weyl, either to the limit-point or to the limit-circle case [3]: $H$ is in the limit-circle case at 0 (respectively at $\infty$ ) if all solutions $y$ of

$$
\begin{equation*}
H y=0 \tag{2.4}
\end{equation*}
$$

are square integrable at 0 (respectively at $\infty$ ), otherwise $H$ is in the limit-point case. According to a simple criterion ([3], p. 231), $H$ is always in the limit-point case at infinity. Since for $x \rightarrow 0$ the two fundamental solutions $y_{1}, y_{2}$ of (2.4) behave as

$$
\begin{equation*}
y_{1}(x) \sim x^{\lambda}, \quad y_{2}(x) \sim x^{-\lambda} \tag{2.5}
\end{equation*}
$$

$H$ is in the limit-point case at 0 for $\lambda \geqslant \frac{1}{2}$ and in the limit-circle case for $\lambda<\frac{1}{2}$ (the case $\lambda=\frac{1}{2}$ is just the superradiant case previously discussed [5]). According to the general theory [3, 4], there exists a unique self-adjoint extension of $H$ in the limitpoint case $\lambda \geqslant \frac{1}{2}$.

In the limit-circle case there are many self-adjoint extensions specified by a boundary condition at $x=0$ of the form ([4], p. 187)

$$
\begin{equation*}
\lim _{x \rightarrow 0}\left\{y(x) \lambda\left(\alpha x^{\lambda}-\beta x^{-\lambda}\right)-y^{\prime}(x) x\left(\alpha x^{\lambda}+\beta x^{-\lambda}\right)\right\}=0 \tag{2.6}
\end{equation*}
$$

We know, however, from reference [1] that the physically important solution $y(x, t)$ corresponding to the quantum optical problem behaves as $x^{\lambda}$ for $x \rightarrow 0$. That requires $\beta=0$ in (2.6). We therefore take for $\lambda<\frac{1}{2}$ the self-adjoint extension defined by

$$
\begin{equation*}
\lim _{x \rightarrow 0}\left\{\lambda x^{\lambda} y(x)-x^{\lambda+1} y^{\prime}(x)\right\}=0 \tag{2.7}
\end{equation*}
$$

In both cases, a self-adjoint Hamiltonian (denoted also by $H$ ) has been obtained. This operator $H$ then generates a unitary group which gives the unique global solution $y(x, t)$ of the Cauchy problem for the Schrödinger equation (2.1) in $L^{2}(0, \infty)$.

Let us now turn to classical mechanics. The Hamiltonian function $H(p, x)$ corresponding to (2.3) is

$$
\begin{equation*}
H(p, x)=x p^{2}+V(x) \tag{2.8}
\end{equation*}
$$

and the corresponding Hamiltonian-Jacobi equation reads

$$
\begin{equation*}
\frac{\partial S}{\partial t}+x\left(\frac{\partial S}{\partial x}\right)^{2}+V(x)=0 \tag{2.9}
\end{equation*}
$$

A complete integral $S$ of this equation depending on two constants of integration $x_{0}$ and $\varepsilon$ is given by the following elliptic integral

$$
\begin{equation*}
\tilde{S}\left(x, x_{0}, \varepsilon ; t, 0\right)=\int_{x_{0}}^{x} d x^{\prime}\left(\eta-\frac{1}{4} x^{\prime 2}-\frac{\varepsilon}{x^{\prime}}-\frac{\lambda^{2}}{x^{\prime 2}}\right)^{1 / 2}+\varepsilon t . \tag{2.10}
\end{equation*}
$$

Then the Cauchy problem of determining a solution $S(x, t)$ of equation (2.9) which satisfies the initial condition

$$
\begin{equation*}
S(x, 0)=S_{0}(x) \tag{2.11}
\end{equation*}
$$

at $t=0$, is solved by forming envelopes of solutions (2.10)

$$
\begin{equation*}
S(x, t)=\tilde{S}\left(x, x_{0}, \varepsilon ; t, 0\right)+S_{0}\left(x_{0}\right), \tag{2.12}
\end{equation*}
$$

where the functions $x_{0}=x_{0}(x, t)$ and $\varepsilon=\varepsilon(x, t)$ are implicitly given by the two relations

$$
\begin{align*}
& \frac{\partial \tilde{S}}{\partial \varepsilon}=0  \tag{2.13}\\
& \frac{\partial \tilde{S}}{\partial x_{0}}+\frac{\partial S_{0}}{\partial x_{0}}=0 \tag{2.14}
\end{align*}
$$

The first relation (2.13) leads to

$$
\begin{equation*}
t=\int_{x_{0}}^{x} R\left(x^{\prime}\right)^{-1 / 2} d x^{\prime} \tag{2.15}
\end{equation*}
$$

with

$$
\begin{equation*}
R(x)=-x^{4}+4 \eta x^{2}-4 \varepsilon x-4 \lambda^{2} . \tag{2.16}
\end{equation*}
$$

This equation (2.15) can be inverted [5], yielding $x_{0}$ (or $x$ respectively)

$$
\begin{equation*}
x_{0}=x+\frac{P^{\prime}(t) R(x)^{1 / 2}+\left(-2 x^{3}+4 \eta x-2 \varepsilon\right)\left(P+\frac{1}{2} x^{2}\right)-x R(x)}{2\left(P+\frac{1}{2} x^{2}\right)^{2}+\frac{1}{2} R(x)} \tag{2.17}
\end{equation*}
$$

Here

$$
\begin{equation*}
P(t)=p(t)-\frac{1}{3} \eta, \quad P^{\prime}(t)=\frac{d p(t)}{d t} \tag{2.18}
\end{equation*}
$$

where

$$
p=p\left(t ; g_{2}, g_{3}\right)
$$

is Weierstrass' $p$-function with invariants

$$
\begin{align*}
& g_{2}=\frac{4}{3} \eta^{2}+4 \lambda^{2} \\
& g_{3}=-\frac{8}{27} \eta^{3}+\frac{8}{3} \eta \lambda^{2}+\varepsilon^{2} . \tag{2.19}
\end{align*}
$$

The equation for $x$ is obtained from (2.17) by the changes $x \leftrightarrow x_{0}, t \leftrightarrow-t$. These equations give the spatial trajectories leading from the point $x$ to the point $x_{0}$ in time $t$ (or reversed) for a given energy $-\varepsilon$.

In order to discuss the second relation (2.14), let us consider the following initial condition (2.11)

$$
\begin{equation*}
S_{0}(x)=-i\left(\lambda \log x-\frac{1}{4} x^{2}\right) \tag{2.20}
\end{equation*}
$$

which will be of interest in the following. Since this $S_{0}$ is no longer real, the corresponding action $S(x, t)$ will be complex. The corresponding complex trajectories do not describe classical motions in the ordinary sense. However, if the quantum mechanical (time dependent) wave function is considered in a classically forbidden region, then the relevant classical trajectories are necessarily complex. Therefore, in a general semi-classical theory, all (real and complex) trajectories have to be used. This point is missed in the book by Maslov [2] (apart from a remark on the last page, 328), which restricts the usefulness of the general apparatus developed there. Inserting (2.20) into equation (2.14), we get the simple expression

$$
\begin{equation*}
\varepsilon=(\eta-\lambda) x_{0} \tag{2.21}
\end{equation*}
$$

This has to be substituted for $\varepsilon$ in (2.17) yielding a transcendental equation for $x_{0}=x_{0}(x, t)$.

Let us consider this equation in some detail for the special value $x=0$, which is of particular interest in the quantum optical problem [1]. In this case the equation for $x_{0}(0, t)$ becomes

$$
\begin{equation*}
x_{0}=-\frac{i \lambda P^{\prime}(t)}{P^{2}-\lambda^{2}+(\eta-\lambda) P} \tag{2.22}
\end{equation*}
$$

We emphasize that $x_{0}$ enters also in the invariant $g_{3}(2.19)$ of the $p$-function through $\varepsilon$ (2.21). The $p$-function and its derivative (2.18) are periodic functions of $t$ with real half-periods $\omega=\omega\left(g_{2}, g_{3}\right)$ depending on the invariants $g_{2}, g_{3}$. Consequently, if $x_{0}\left(t_{1}\right)$ is a solution of (2.22) for $t=t_{1}$, it is also a solution for $t=t_{1}+\omega$. But the periods vary with the value of $x_{0}$. Therefore $x_{0}(t)$ is not periodic in $t$, but it can be constructed from the values in the interval $0 \leqslant t \leqslant \omega_{0}=\omega\left(x_{0}=0\right)$ by a nonuniform shear transformation (see Figure 1). For small times, $x_{0}(t)$ is unique and shows an oscillatory behavior. Since the period $\omega\left(x_{0}\right)$ increases with increasing $\left|x_{0}\right|$, the larger values of $\left|x_{0}\right|$ are shifted more to the right than the smaller ones during the next oscillation. This leads to a fall over of the wave, which occurs already in the second oscillation, in general. In Figure 1 numerical values of $x_{0}(0, t)$ are shown for $\lambda=\frac{1}{2}, \eta=\frac{3}{2}$. We see that $x_{0}(t)$ necessarily becomes multivalued for larger $t$. The branch-points where new branches $x_{0 k}(t)$ appear (or disappear), i.e. where

$$
\begin{equation*}
\frac{\partial x_{0}}{\partial t}= \pm \infty \tag{2.23}
\end{equation*}
$$

are the focal points. The number of branches increases linearly with $t$.
In the general case $x \neq 0$ the situation is qualitatively the same. Since $x_{0}(x, t)$ is multivalued, $S(x, t)(2.12)$ becomes multivalued as well. We therefore get finitely many different solutions $S_{k}(x, t)$ of the Cauchy problem for the Hamilton-Jacobi equation. But there exists a finite time interval [ $0, T(x)$ ) (larger than one oscillation period) where the solution $x_{0}(x, t)$ is unique. Moreover, the focal points $T(x)$ are


Figure 1:
$x_{0}(x, t)$ for $x=0, \lambda=0.5, \eta=1.5$ as a function of $t$. The arrows illustrate the shear transformation. The focal points are marked by little circles.
bounded away from 0

$$
\begin{equation*}
T_{0}=\inf _{0 \leqslant x<\infty} T(x)>0 \tag{2.24}
\end{equation*}
$$

As in the case of ordinary mechanical systems we could try to solve equation (2.15) for $\varepsilon$

$$
\begin{equation*}
\varepsilon=\varepsilon\left(x, x_{0}, t\right) \tag{2.25}
\end{equation*}
$$

in order to get the action $\tilde{S}(2.10)$ as a function of the coordinates $x_{0}, x$ and $t$. Equation (2.15) can indeed be inverted in the form ([5] equation 4.19)

$$
\begin{equation*}
\varepsilon=\left(x+x_{0}\right)\left(\frac{1}{2} x x_{0}-P\right)+Q^{1 / 2} \tag{2.26}
\end{equation*}
$$

with

$$
\begin{equation*}
Q=\left[\left(x+x_{0}\right)^{2}-4 \eta-4 P\right]\left[\frac{1}{4} x^{2} x_{0}^{2}+\lambda^{2}-P x x_{0}\right] . \tag{2.27}
\end{equation*}
$$

But this transcendental equation for $\varepsilon$ has infinitely many solutions, in general, even for arbitrarily small $t$ ! This can be seen as follows: For fixed $x, x_{0}, \varepsilon\left(x, x_{0}, t\right)$ is a monotonely decreasing function of $t(\varepsilon(t=0)=+\infty)$, following one branch. In the same manner as was discussed above, infinitely many new branches can be generated from a given one by non-uniform shear transformations. The crucial point is that the real period $\omega$ of the shear translation goes to 0 for $t \rightarrow 0$, which can be deduced from the expansion of the $p$-function for small $t$. Consequently, the branches approach each other for $t \rightarrow 0$. Then there are infinitely many different values of $\varepsilon\left(x, x_{0}, t\right)$ for all $t>0$. In Figure 2, $\varepsilon\left(x, x_{0}, t\right)$ as a function of $t$ is depicted for fixed $x, x_{0}$ in a typical case.


Figure 2:
Energy $\varepsilon\left(0, x_{0}, t\right)$ for $x_{0}=0.4 i, \lambda=0.25, \eta=1.25$ as a function of $t$. The arrows illustrate the shear transformation.

This behavior of the energy $\varepsilon$ has some unusual consequences. Given two points $x_{0}$ and $x$ in space, then there exist infinitely many trajectories leading from $x_{0}$ to $x$ in a given time $t$, and this remains true for arbitrarily small $t$. That is in sharp contrast to the behavior of a ordinary point mass (with kinetic energy $p^{2} / 2 m$ ). At least if the potential $V(x)$ is bounded, there is only one trajectory from $x_{0}$ to $x$ for sufficiently small $t$ ([7], theorem 3.1). At the same time the action $\widetilde{S}_{k}\left(x, x_{0}, t\right)$ has infinitely many different branches for all $t>0$. Then, it is quite complicated to derive a semi-classical approximation for the Green's function $K\left(x, x_{0}, \mathrm{t}\right)$ of the Schrödinger equation (2.1) even for small $t$, as can be done in ordinary quantum mechanics [7]. Therefore, it seems to be rather difficult to convert the formal arguments of reference [1] into a proof. For this reason, we proceed a somewhat different way in the following sections.

## 3. The local semi-classical solution

It is our aim to construct an asymptotic approximation for $h \rightarrow 0$ to the solution $y(x, t)$ of the Cauchy problem for the Schrödinger equation (2.1). Following Maslov [6], we consider the initial condition

$$
\begin{equation*}
y(x, 0)=y_{0}(x)=\rho(x) \exp \frac{1}{h} \varphi(x) \tag{3.1}
\end{equation*}
$$

In this section, we assume $\rho$ to be in $L^{2}(0, \infty)$ and $\operatorname{Re} \varphi(x)$ bounded from above, in addition $\rho$ and $\varphi$ are supposed to have analytic continuations into the complex $x$-plane. For the physically most interesting initial state (the so-called fully excited state of quantum optics) we have [1]

$$
\begin{equation*}
y_{0}(x)=\exp \frac{1}{h}\left(\lambda \log x-\frac{1}{4} x^{2}\right), \tag{3.2}
\end{equation*}
$$

which is not precisely of the form (3.1). But if $y_{0}$ is normalized such that the $L^{2}(0, \infty)-$ norm is 1

$$
\begin{equation*}
\left\|y_{0}\right\|^{2}=1 \tag{3.3}
\end{equation*}
$$

then in the limit $h \rightarrow 0, y_{0}^{2}(x)$ goes to

$$
\begin{equation*}
y_{0}^{2}(x) \rightarrow \delta(x-\sqrt{2 \lambda}) \tag{3.4}
\end{equation*}
$$

in the sense of distributions. Therefore, if we multiply by a continuous $\rho(x) \in L^{2}(0, \infty)$ with $\rho(\sqrt{2 \lambda})=1$, the error is exponentially small for $h \rightarrow 0$.

For the solution $y(x, t)$ of the Cauchy problem, we now make the following ansatz

$$
\begin{equation*}
\mathrm{y}(x, t)=u(x, t) \exp \frac{i}{h} S(x, t) \tag{3.5}
\end{equation*}
$$

where $S(x, t)$ is the solution of the Hamilton-Jacobi equation (2.9) satisfying the initial condition

$$
\begin{equation*}
S(x, 0)=S_{0}(x)=-i \varphi(x) . \tag{3.6}
\end{equation*}
$$

As was pointed out in the last section, $S(x, t)$ is uniquely determined for small $t<T_{0}$ (2.24).

In particular, for the fully excited initial state (3.2) this has been investigated (2.20). If (3.5) is substituted into the Schrödinger equation (2.1), one gets the following equation for $u(x, t)$

$$
\begin{align*}
& i h \frac{\partial u}{\partial t}=\left(\frac{\partial S}{\partial t}+x\left(\frac{\partial S}{\partial x}\right)^{2}+V(x)\right) u  \tag{3.7}\\
& -i h\left(x \frac{\partial^{2} S}{\partial x^{2}} u+2 x \frac{\partial S}{\partial x} \frac{\partial u}{\partial x}+\frac{\partial S}{\partial x} u\right)-h^{2}\left(x \frac{\partial^{2} u}{\partial x^{2}}+\frac{\partial u}{\partial x}\right)
\end{align*}
$$

where the first term on the right side vanishes due to the choice of $S(x, t)$. The corresponding initial condition is

$$
\begin{equation*}
u(x, 0)=\rho(x) . \tag{3.8}
\end{equation*}
$$

Let us consider the function $x_{0}(x, t)$ which is implicitly given in terms of the trajectory

$$
\begin{equation*}
x=x\left(x_{0}, t\right) . \tag{3.9}
\end{equation*}
$$

We calculate the time derivative of $x_{0}$ (with $x$ fixed) from (3.9)

$$
\begin{equation*}
\frac{\partial x_{0}}{\partial t}=-\frac{\frac{\partial x}{\partial t}}{\frac{\partial x}{\partial x_{0}}} \tag{3.10}
\end{equation*}
$$

and use

$$
\begin{equation*}
\frac{\partial x}{\partial t}=\frac{\partial H}{\partial p}=2 x p=2 x \frac{\partial S}{\partial x} \tag{3.11}
\end{equation*}
$$

obtaining

$$
\begin{equation*}
\frac{\partial x_{0}}{\partial t}=-\frac{\partial x_{0}}{\partial x} 2 x \frac{\partial S}{\partial x} \tag{3.12}
\end{equation*}
$$

Differentiating this with respect to $x$, we get the relation

$$
\begin{equation*}
\frac{\partial}{\partial x} \frac{\partial x_{0}}{\partial t}=\frac{\partial}{\partial t} \frac{\partial x_{0}}{\partial x}=-\frac{\partial^{2} x_{0}}{\partial x^{2}} 2 x \frac{\partial S}{\partial x}-2 \frac{\partial x_{0}}{\partial x} \frac{\partial S}{\partial x}-2 x \frac{\partial x_{0}}{\partial x} \frac{\partial^{2} S}{\partial x^{2}} \tag{3.13}
\end{equation*}
$$

For the solution $u(x, t)$ of (3.7), (3.8), we now make the ansatz

$$
\begin{equation*}
u(x, t)=\left(\frac{\partial x_{0}}{\partial x}\right)^{1 / 2} v_{h}(x, t) \tag{3.14}
\end{equation*}
$$

Substituting into equation (3.7) and using (3.13), we find the following equation for $v_{h}(x, t)$

$$
\begin{equation*}
\frac{\partial v_{h}}{\partial t}+2 x \frac{\partial S}{\partial x} \frac{\partial v_{h}}{\partial x}=i h\left(\frac{\partial x_{0}}{\partial x}\right)^{-1 / 2} \frac{\partial}{\partial x} x \frac{\partial}{\partial x}\left(\frac{\partial x_{0}}{\partial x}\right)^{1 / 2} v_{h} \tag{3.15}
\end{equation*}
$$

The initial condition now reads

$$
\begin{equation*}
v_{h}(x, 0)=\rho(x) \tag{3.16}
\end{equation*}
$$

For the left side of equation (3.15) we can write

$$
\begin{equation*}
\frac{\partial v_{h}}{\partial t}+\frac{\partial x}{\partial t} \frac{\partial v_{h}}{\partial x}=\frac{d}{d t} v_{h}\left(x\left(x_{0}, t\right), t\right) \tag{3.17}
\end{equation*}
$$

which is the time derivative along the trajectory (3.9). Then (3.15) and (3.16) are equivalent to the integral equation

$$
\begin{equation*}
v_{h}(x, t)=\rho\left(x_{0}\right)+i h \int_{0}^{t} d t_{1}\left(\frac{\partial x_{0}}{\partial x}\right)^{-1 / 2} \frac{\partial}{\partial x} x \frac{\partial}{\partial x}\left(\frac{\partial x_{0}}{\partial x}\right)^{1 / 2} v_{h}\left(x, t_{1}\right) \tag{3.18}
\end{equation*}
$$

where the integral is taken along the trajectory (3.9), that means $x$ has to be substituted by $x\left(x_{0}, t_{1}\right)$. Since for $0 \leqslant t<T_{0}$ there exists a unique trajectory, it follows that

$$
\begin{equation*}
\frac{\partial x}{\partial x_{0}} \neq 0 \tag{3.19}
\end{equation*}
$$

Therefore, the factors appearing in (3.14) and (3.18) are bounded. If the time $t$ approaches a focal point, then according to (2.23) and (3.10)

$$
\begin{equation*}
\frac{\partial x}{\partial x_{0}} \rightarrow 0 \tag{3.20}
\end{equation*}
$$

which is usually taken as the definition of a focal point [7].
We must now control the $h$-dependence of $v_{h}(x, t)$. For this purpose, let us perform a scaling transformation

$$
\begin{array}{ll}
x^{\prime}=x h^{-1 / 2}, & t^{\prime}=t h^{1 / 2} \\
\lambda^{\prime}=\lambda / h, & \eta^{\prime}=\eta / h \tag{3.21}
\end{array}
$$

in equation (3.15), which leads to

$$
\begin{equation*}
\frac{\partial v}{\partial t^{\prime}}+\frac{\partial x^{\prime}}{\partial t^{\prime}} \frac{\partial v}{\partial x^{\prime}}=i\left(\frac{\partial x_{0}^{\prime}}{\partial x^{\prime}}\right)^{-1 / 2} \frac{\partial}{\partial x^{\prime}} x^{\prime} \frac{\partial}{\partial x^{\prime}}\left(\frac{\hat{c} x_{0}^{\prime}}{\partial x^{\prime}}\right)^{1 / 2} v . \tag{3.22}
\end{equation*}
$$

Since this equation no longer contains the parameter $h$ and the initial condition (3.16) is independent of $h$ too, it follows that

$$
\begin{equation*}
v_{h}(x, t)=v\left(x h^{-1 / 2}, t h^{1 / 2}\right) . \tag{3.23}
\end{equation*}
$$

Then, since $v\left(x^{\prime}, t^{\prime}\right)$ is bounded for $x^{\prime} \rightarrow \infty, t^{\prime} \rightarrow 0, v_{h}(x, t)$ is bounded with respect to $h$ uniformly in $x$ and $0 \leqslant t \leqslant t_{1}<T_{0}$, together with its first and second $x$-derivatives. Consequently, we may conclude from the integral equation (3.18) that

$$
\begin{equation*}
v_{h}(x, t)=\rho\left(x_{0}(x, t)\right)+0(h) \tag{3.24}
\end{equation*}
$$

for $h \rightarrow 0$ and $x$ fixed.
Summing up, we have shown that

$$
\begin{equation*}
y(x, t)=\left(\frac{\partial x_{0}}{\partial x}\right)^{1 / 2} e^{i /(h) S(x, t)}\left[\rho\left(x_{0}\right)+0(h)\right] \tag{3.25}
\end{equation*}
$$

for $t<T_{0}$. Since the error term is square integrable, we have also an approximation in the $L^{2}$ sense

$$
\begin{equation*}
\left\|y(x, t)-y_{c}(x, t)\right\|=0(h) \tag{3.26}
\end{equation*}
$$

where

$$
\begin{equation*}
y_{c}(x, t)=\rho\left(x_{0}\right)\left(\frac{\partial x_{0}}{\partial x}\right)^{1 / 2} \exp \left(\frac{i}{h} S(x, \mathrm{t})\right) \tag{3.27}
\end{equation*}
$$

is the local semi-classical solution.

## 4. The global semi-classical solution

The local semi-classical solution (3.27) diverges at the first focal point according to (3.20). However, beyond the focal points the expression (3.27) is again finite but not unique. The exact solution $y(x, t)$, on the other hand, remains regular at the focal point. There rises the connection problem of continuing $y_{c}(x, t)$ in $t$ through the focal points in such a way that it approximates $y(x, t)$ beyond the focal points.

For technical reasons, we now assume that $\rho(x)$ and $\varphi(x)$ appearing in the initial condition (3.1) are analytic functions of $x$, regular in a certain region $\Omega^{\prime}$ of the complex $x$-plane which contains a part of the positive real axis and will be further specified below. Outside $\Omega^{\prime} \rho(x)$ is assumed to be $C^{2}$ with compact support $\Omega \supset \Omega^{\prime}$, while $\varphi(x)$ is $C^{2}$ in $\Omega \backslash \Omega^{\prime}$. We remark that these assumptions are not too restrictive, because a $\rho(x)$ with non-compact support can be decomposed by means of a decomposition of unity

$$
1=\sum_{k} \rho_{k}(x), \quad \rho_{k} \in C_{0}^{2}
$$

The focal points $T(x)$ of the local semi-classical solution (3.27) lie in the time interval $\left[T_{1}, T_{2}\right]$

$$
\begin{align*}
T_{1} & =\min _{x_{0} \in \Omega} T(x) \\
T_{2} & =\max _{x_{0} \in \Omega} T(x) \tag{4.1}
\end{align*}
$$

Following a method of Maslov ([2], p. 285), we bridge the gap between $T_{1}$ and $T_{2}$ by going over to the momentum space. Let us choose $t_{1}<T_{1}$, so that the semiclassical solution $y_{c}\left(x, t_{1}\right)(3.27)$ is still valid at $t_{1}$. We compute the Fourier transform

$$
\begin{equation*}
\hat{y}_{c}\left(p, t_{1}\right)=(2 \pi h)^{-1 / 2} \int d x \mathrm{e}^{-i /(h) p x}\left(\frac{\partial x_{0}}{\partial x}\right)^{1 / 2} \mathrm{e}^{i /(h) S\left(x, t_{1}\right)} \rho\left(x_{0}\right), \tag{4.2}
\end{equation*}
$$

where the integral actually runs only over the compact set $\Omega_{x}$

$$
\begin{equation*}
\Omega_{x}=\left\{x\left(x_{0}, t_{1}\right) \in \mathbb{R}_{+}^{1} \mid x_{0} \in \Omega\right\} . \tag{4.3}
\end{equation*}
$$

The integrand in (4.2) is analytic in

$$
\begin{equation*}
\Omega_{x}^{\prime}=\left\{x\left(x_{0}, t_{1}\right) \mid x_{0} \in \Omega^{\prime}\right\} . \tag{4.4}
\end{equation*}
$$

The integral (4.2) will be evaluated by the saddle-point method. A saddle-point $x_{1}$ is a solution of the equation

$$
\begin{equation*}
\frac{\partial S\left(x, t_{1}\right)}{\partial x}=p \tag{4.5}
\end{equation*}
$$

that means $p$ must be the momentum at time $t_{1}$ corresponding to the unique trajectory (3.9)

$$
\begin{align*}
& p=p\left(x_{0}, t_{1}\right)  \tag{4.6}\\
& x_{1}=x\left(x_{0}, t_{1}\right)
\end{align*}
$$

Let $\Lambda$ be the set of momentum values

$$
\begin{equation*}
\Lambda=\left\{p\left(x_{0}, t_{1}\right) \mid x\left(x_{0}, t_{1}\right) \in \Omega_{x}\right\} \tag{4.7}
\end{equation*}
$$

Then, for $p \in \Lambda$ there exists a unique saddle-point $x_{1}$, generally complex. For $p \in \Lambda^{\prime}$

$$
\begin{equation*}
\Lambda^{\prime}=\left\{p\left(x_{0}, t_{1}\right) \in \Lambda \mid x_{0} \in \Omega^{\prime}\right\} \tag{4.8}
\end{equation*}
$$

the saddle-point $x_{1}$ lies in the analyticity domain $\Omega_{x}^{\prime}$ (4.4). The path of integration can then be deformed inside $\Omega_{x}^{\prime}$ so that it runs through $x_{1}$. Now the integral inside $\Omega_{x}^{\prime}$ can be evaluated by the complex saddle-point method ([8], p. 232) up to an exponentially small error, the integral in the rest of $\Omega_{x}$ is $0\left(h^{2}\right)$ because $\rho \in C^{2}$. This yields

$$
\begin{align*}
\hat{y}_{c}\left(p, t_{1}\right)= & \left(\frac{2 \pi h}{\frac{\partial^{2}}{\partial x^{2}} S\left(x, t_{1}\right)}\right)_{x_{1}}^{1 / 2}(2 \pi h)^{-1 / 2} \mathrm{e}^{-i /(h) p x_{1}}\left(\frac{\partial x_{0}}{\partial x}\right)_{x_{1}}^{1 / 2} \\
& \left.\times \mathrm{e}^{i /(h) S\left(x_{1}, t_{1}\right.}\right) \rho\left(x_{0}^{\prime}\right)+0\left(h^{3 / 2}\right), \tag{4.9}
\end{align*}
$$

with $x_{0}^{\prime}=x_{0}\left(x_{1}, t_{1}\right)$. The error term in (4.9) can be estimated in the sup-norm or in the $L^{2}$ norm. In the following, all estimates are understood in the $L^{2}$ sense. For $p \notin \Lambda$ there exists no saddle point, therefore

$$
\begin{equation*}
\hat{y}_{c}\left(p, t_{1}\right)=0\left(h^{3 / 2}\right) . \tag{4.10}
\end{equation*}
$$

The contribution from $p \in \Lambda \backslash \Lambda^{\prime}$ can be made arbitrarily small in the $L^{2}$ sense, say
$0\left(h^{3 / 2}\right)$ also, by choosing $\Omega \backslash \Omega^{\prime}$ small enough. The error $0\left(h^{3 / 2}\right)$ can be neglected compared with the error $0(h)$ of the local semi-classical solution (3.27) (see (4.22) below). From (4.6) we get

$$
\begin{equation*}
x_{0}=x_{0}\left(p, t_{1}\right) \tag{4.11}
\end{equation*}
$$

and

$$
\begin{equation*}
x_{1}=x_{1}\left(p, t_{1}\right) \tag{4.12}
\end{equation*}
$$

Using (4.5) and the chain rule, we can write (4.9) as

$$
\begin{equation*}
\hat{y}_{c}\left(p, t_{1}\right)=\left(\frac{\partial x_{0}}{\partial p}\right)^{1 / 2} \mathrm{e}^{i /(h) \boldsymbol{S}\left(p, t_{1}\right)} \rho\left(x_{0}^{\prime}\right)+0\left(h^{3 / 2}\right) \tag{4.13}
\end{equation*}
$$

where

$$
\begin{equation*}
S\left(p, t_{1}\right)=S\left(x_{1}\left(p, t_{1}\right), t_{1}\right)-p x_{1}\left(p, t_{1}\right) \tag{4.14}
\end{equation*}
$$

According to (4.5), this is a Legendre transformation. We have

$$
\begin{equation*}
\frac{\partial S\left(p, t_{1}\right)}{\partial p}=-x_{1} \tag{4.15}
\end{equation*}
$$

Let $\Lambda_{p}$ be the compact set on the real $p$-axis which is in $\Lambda^{\prime}(4.8)$, and $\chi_{p}(p) \in C_{0}^{\infty}\left(\mathbb{R}^{1}\right)$

$$
\chi_{p}(p)=\left\{\begin{array}{rl}
1 & p \in \Lambda_{p}  \tag{4.16}\\
<1 & p \notin \Lambda_{p} \\
0 & p \notin \operatorname{supp} \chi_{p} \supset \Lambda_{p}
\end{array}\right.
$$

and similarly

$$
\chi_{x}(x)=\left\{\begin{array}{rl}
1 & x \in \Omega_{x}  \tag{4.17}\\
<1 & x \notin \Omega_{x} \\
0 & x \notin \operatorname{supp} \chi_{x} \supset \Omega_{x}
\end{array}\right.
$$

In addition, we define $\chi_{x}$ and $\chi_{p}=1$ in some complex region of the $p$ or $x$ plane, respectively, specified below. The function

$$
\begin{equation*}
\tilde{y}_{c}\left(p, t_{1}\right)=\chi_{p}(p)\left(\frac{\partial x_{0}}{\partial p}\right)^{1 / 2} \mathrm{e}^{i /(h) S\left(p, t_{1}\right)} \rho\left(x_{0}^{\prime}\right) \tag{4.18}
\end{equation*}
$$

is approximately equal to $\hat{y}_{c}\left(p, t_{1}\right)$ in $L_{2}\left(\mathbb{R}^{1}\right)$ up to $0\left(h^{3 / 2}\right)$. It then follows from the Fourier inversion theorem in $L^{2}$ that

$$
\begin{equation*}
y_{c}\left(x, t_{1}\right)=(2 \pi h)^{-1 / 2} \int d p \tilde{y}_{c}\left(p, t_{1}\right) \mathrm{e}^{i /(h) p x}+0\left(h^{3 / 2}\right) \tag{4.19}
\end{equation*}
$$

Multiplying with $\chi_{x}$ (4.17), we have also

$$
\begin{equation*}
y_{c}\left(x, t_{1}\right)=\chi_{x}(x)(2 \pi h)^{-1 / 2} \int d p \tilde{y}_{c}\left(p, t_{1}\right) \mathrm{e}^{i /(h) p x}+0\left(h^{3 / 2}\right) \tag{4.20}
\end{equation*}
$$

We now consider $t>t_{1}$. Let $y_{p}(x, t)$ be the solution of the Schrödinger equation (2.1) with the initial condition

$$
\begin{equation*}
y_{p}\left(x, t_{1}\right)=\chi_{x}(x) \mathrm{e}^{i /(h) p x} \tag{4.21}
\end{equation*}
$$

at $t=t_{1}$. It follows from (3.27) and (4.20) that

$$
\begin{equation*}
y(x, t)=(2 \pi h)^{-1 / 2} \int d p \tilde{y}_{c}\left(p, t_{1}\right) y_{p}(x, \mathrm{t})+0(h) \tag{4.22}
\end{equation*}
$$

for arbitrary $t \geqslant t_{1}$, due to the unitarity of the exact time evolution operator. Since $y_{p}\left(x, t_{1}\right)$ is of the form (3.1), we can apply the results of Section 3 and compute its local semi-classical approximation $y_{p c}(x, t)$. We have to consider the trajectory

$$
\begin{align*}
& x^{\prime}=x^{\prime}\left(x_{2}, p, t^{\prime}\right)  \tag{4.23}\\
& p^{\prime}=p^{\prime}\left(x_{2}, p, t^{\prime}\right), t_{1} \leqslant t^{\prime} \leqslant t \tag{4.24}
\end{align*}
$$

with $x^{\prime}=x_{2}, p^{\prime}=p$ at $t=t_{1}$. For every given $x=x(t)$ and $p=p^{\prime}\left(t_{1}\right) \in \operatorname{supp} \chi_{p}$ there exists a unique initial coordinate $x_{2}$

$$
\begin{equation*}
x_{2}=x_{2}(x, p, t) \tag{4.25}
\end{equation*}
$$

if $t$ is not too big, say

$$
\begin{equation*}
t-t_{1} \leqslant \Delta t \tag{4.26}
\end{equation*}
$$

Then there also exists a unique action $S_{p}(x, t)$ given by (2.12)

$$
\begin{equation*}
S_{p}(x, t)=\widetilde{S}\left(x, x_{2} ; t, t_{1}\right)+p x_{2} \tag{4.27}
\end{equation*}
$$

and a corresponding local semi-classical solution

$$
\begin{equation*}
y_{p c}(x, t)=\chi_{x}\left(x_{2}\right)\left(\frac{\partial x_{2}}{\partial x}\right)^{1 / 2} \exp \left[\frac{i}{h} S_{p}(x, t)\right] \tag{4.28}
\end{equation*}
$$

The energy $\varepsilon$ has been dropped in the argument of $\tilde{S}$, because it is now explicitly given by

$$
\begin{equation*}
-\varepsilon=x_{2} p^{2}+V\left(x_{2}\right) \tag{4.29}
\end{equation*}
$$

Using

$$
\begin{equation*}
\frac{\partial \widetilde{S}}{\partial x_{2}}=-p \tag{4.30}
\end{equation*}
$$

it follows from (4.27) that

$$
\begin{equation*}
\frac{\partial S_{p}}{\partial p}=\left(\frac{\partial \tilde{S}}{\partial x_{2}}+p\right) \frac{\partial x_{2}}{\partial p}+x_{2}=x_{2} \tag{4.31}
\end{equation*}
$$

Finally, substituting $y_{p c}$ (4.28) for $y_{p}$ into (4.22), we obtain with (4.18)

$$
\begin{align*}
y(x, t)= & (2 \pi h)^{-1 / 2} \int d p \chi_{p}(p)\left(\frac{\partial x_{0}}{\partial p}\right)^{1 / 2} \mathrm{e}^{i /(h) S\left(p, t_{1}\right)} \rho\left(x_{0}^{\prime}\right) \\
& \times \chi_{x}\left(x_{2}\right)\left(\frac{\partial x_{2}}{\partial x}\right)^{1 / 2} \mathrm{e}^{i /(h) S_{p}(x, t)}+0(h) \tag{4.32}
\end{align*}
$$

This integral will again be computed by the saddle-point method. The saddle-points are solutions of

$$
\begin{equation*}
\frac{\partial S\left(p, t_{1}\right)}{\partial p}+\frac{\partial S_{p}(x, t)}{\partial p}=0 \tag{4.33}
\end{equation*}
$$

or with (4.15) and (4.31)

$$
\begin{equation*}
x_{1}\left(p, t_{1}\right)=x_{2}(x, p) \tag{4.34}
\end{equation*}
$$

Consequently, $p$ must be chosen in such a way that the two trajectories (3.9) from time 0 to $t_{1}$ and (4.23) from $t_{1}$ to $t$ join together, forming one trajectory from 0 to $t$. Since the latter is no longer unique for $t>T_{1}$, there exists now more than one saddle point (the argument given in Maslov's book [2] p. 293 is wrong at this point). The exponent in (4.32) simplifies at the saddle-points $p_{k}$ :

$$
\begin{align*}
S\left(p_{k}, t_{1}\right)+S_{p_{k}}(x, t)= & -p_{k} x_{1}\left(p_{k}\right)+\widetilde{S}\left(x_{1}, x_{0} ; t_{1}, 0\right) \\
& +\stackrel{S}{0}\left(x_{0}\right)+\widetilde{S}\left(x, x_{2} ; t, t_{1}\right)+p_{k} x_{2}\left(p_{k}\right) \\
= & \stackrel{\tilde{S}_{k}\left(x, x_{0}\left(p_{k}\right) ; t, 0\right)+S_{0}\left(x_{0}\left(p_{k}\right)\right)}{=} S_{k}(x, t) .
\end{align*}
$$

The domains of analyticity for $\chi_{p}$ and $\chi_{x}$ must be chosen big enough such that all (complex) saddle-points $p_{k}$ lie inside. Assuming $t_{1}+\Delta t>T_{2}$, which can always be arranged by a suitable restriction of the support $\Omega$ of $\rho$, we can compute the integral (4.32) for $T_{2}<t<t_{1}+\Delta t$ by the saddle-point method

$$
\begin{align*}
y(x, t)= & \sum_{k}\left(\frac{2 \pi h}{\frac{\partial x_{2}}{\partial p}-\frac{\partial x_{1}}{\partial p}} \frac{1}{2 \pi h} \frac{\partial x_{0}}{\partial p} \frac{\partial x_{2}}{\partial x}\right)_{p_{k}}^{1 / 2} \\
& \times \mathrm{e}^{i /(h) S_{k}(x, t)} \rho\left(x_{0 k}\right)+0(h) \tag{4.36}
\end{align*}
$$

where

$$
\begin{align*}
& x_{0 k}=x_{0}\left(x_{1}\left(p_{k}, t_{1}\right), t_{1}\right)=x_{0}\left(x_{2}\left(x, p_{k}\right), t_{1}\right) \\
& \quad \stackrel{\text { def }}{=} x_{0 k}(x, t) \tag{4.37}
\end{align*}
$$

Here $p$ (respectively $p_{k}$ ) appears as a parameter which must be determined by $x$, so that everything on the right side of (4.36) depends ultimately on $x$ only.

We set $t^{\prime}=t_{1}$ in (4.24) and solve (4.25) for $p$ to obtain

$$
\begin{equation*}
p=p^{\prime}\left(x_{2}, p\left(x, x_{2}, t\right), t_{1}\right) \stackrel{\text { def }}{=} p\left(x_{2}, x\right) \tag{4.38}
\end{equation*}
$$

or because of (4.34)

$$
\begin{equation*}
p-p\left(x_{1}(p), x\right)=0 \tag{4.39}
\end{equation*}
$$

omitting the time arguments which are now of no interest. This equation (4.39) gives $p$ as an implicit function of $x$. The $x$-derivative then is

$$
\begin{equation*}
\frac{\partial p}{\partial x}=\frac{\left(\frac{\partial p}{\partial x}\right)_{x_{2}}}{1-\left(\frac{\partial p}{\partial x_{2}}\right)_{x} \frac{\partial x_{1}}{\partial p}}=\frac{\left(\frac{\partial x_{2}}{\partial p}\right)_{x}\left(\frac{\partial p}{\partial x}\right)_{x_{2}}}{\left(\frac{\partial x_{2}}{\partial p}\right)_{x}-\frac{\partial x_{1}}{\partial p}} \tag{4.40}
\end{equation*}
$$

Using this expression in (4.36), we get the final result ([1], equation (2.31))

$$
\begin{equation*}
y(x, t)=\sum_{k}\left(\frac{\partial x_{0 k}}{\partial x}\right)^{1 / 2} \mathrm{e}^{i /(h) S_{k}(x, t)} \rho\left(x_{0 k}\right)+0(h) . \tag{4.41}
\end{equation*}
$$

This formula shows that the local semi-classical solution can be continued beyond the first focal points in the most plausible way by adding up the contributions
of the different branches. The same arguments can then be used at the following focal points, at which one considers each branch separately. In this way, the global semi-classical solution is established. As we have seen, at each focal point an additional error $0(h)$ comes in. Since the number of focal points increases linearly with $t$, the total error increases in the same way, if all error terms have the same order of magnitude. In the quantum optical application [1] $h$ is proportional to $N^{-1}$ ( $N=$ number of radiating atoms). The global semi-classical solution then breaks down after $0(N)$ focal points or after $0(N)$ oscillations. This is precisely what one finds in numerical calculations for the Dicke model (see also [5] p. 85).

A final remark is concerned with the nature of the semi-classical approximation. The semi-classical solution $y_{c}(x, t)$ approaches $y(x, t)$ for $h \rightarrow 0$ in the $L^{2}$ sense, but in the applications [1] the behavior of $y(x, t)$ in the neighborhood of $x=0$ is required. However, going over to a subsequence $h_{j} \rightarrow 0$, we get convergence almost everywhere, and this is sufficient for the applications.

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